# Complex Analytic and Differential Geometry 

Jean-Pierre Demailly

Université de Grenoble I Institut Fourier, UMR 5582 du CNRS 38402 Saint-Martin d'Hères, France

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## Chapter I

## Complex Differential Calculus and Pseudoconvexity


#### Abstract

This introductive chapter is mainly a review of the basic tools and concepts which will be employed in the rest of the book: differential forms, currents, holomorphic and plurisubharmonic functions, holomorphic convexity and pseudoconvexity. Our study of holomorphic convexity is principally concentrated here on the case of domains in $\mathbb{C}^{n}$. The more powerful machinery needed for the study of general complex varieties (sheaves, positive currents, hermitian differential geometry) will be introduced in Chapters II to V. Although our exposition pretends to be almost self-contained, the reader is assumed to have at least a vague familiarity with a few basic topics, such as differential calculus, measure theory and distributions, holomorphic functions of one complex variable, .... Most of the necessary background can be found in the books of [Rudin 1966] and [Warner 1971]; the basics of distribution theory can be found in Chapter I of [Hörmander 1963]. On the other hand, the reader who has already some knowledge of complex analysis in several variables should probably bypass this chapter.


## § 1. Differential Calculus on Manifolds

## §1.A. Differentiable Manifolds

The notion of manifold is a natural extension of the notion of submanifold defined by a set of equations in $\mathbb{R}^{n}$. However, as already observed by Riemann during the 19th century, it is important to define the notion of a manifold in a flexible way, without necessarily requiring that the underlying topological space is embedded in an affine space. The precise formal definition was first introduced by H. Weyl in [Weyl 1913].

Let $m \in \mathbb{N}$ and $k \in \mathbb{N} \cup\{\infty, \omega\}$. We denote by $\mathscr{C}^{k}$ the class of functions which are $k$-times differentiable with continuous derivatives if $k \neq \omega$, and by $C^{\omega}$ the class of real analytic functions. A differentiable manifold $M$ of real dimension $m$ and of class $\mathscr{C}^{k}$ is a topological space (which we shall always assume Hausdorff and separable, i.e. possessing a countable basis of the topology), equipped with an atlas of class $\mathscr{C}^{k}$ with values in $\mathbb{R}^{m}$. An atlas of class $\mathscr{C}^{k}$ is a collection of homeomorphisms $\tau_{\alpha}: U_{\alpha} \longrightarrow V_{\alpha}, \alpha \in I$, called differentiable charts, such that $\left(U_{\alpha}\right)_{\alpha \in I}$ is an open covering of $M$ and $V_{\alpha}$ an open subset of $\mathbb{R}^{m}$, and such that for all $\alpha, \beta \in I$ the transition map

$$
\begin{equation*}
\tau_{\alpha \beta}=\tau_{\alpha} \circ \tau_{\beta}^{-1}: \tau_{\beta}\left(U_{\alpha} \cap U_{\beta}\right) \longrightarrow \tau_{\alpha}\left(U_{\alpha} \cap U_{\beta}\right) \tag{1.1}
\end{equation*}
$$

is a $C^{k}$ diffeomorphism from an open subset of $V_{\beta}$ onto an open subset of $V_{\alpha}$ (see Fig. 1). Then the components $\tau_{\alpha}(x)=\left(x_{1}^{\alpha}, \ldots, x_{m}^{\alpha}\right)$ are called the local coordinates on $U_{\alpha}$ defined by the chart $\tau_{\alpha}$; they are related by the transition relation $x^{\alpha}=\tau_{\alpha \beta}\left(x^{\beta}\right)$.


Fig. I-1 Charts and transition maps

If $\Omega \subset M$ is open and $s \in \mathbb{N} \cup\{\infty, \omega\}, 0 \leqslant s \leqslant k$, we denote by $C^{s}(\Omega, \mathbb{R})$ the set of functions $f$ of class $C^{s}$ on $\Omega$, i.e. such that $f \circ \tau_{\alpha}^{-1}$ is of class $C^{s}$ on $\tau_{\alpha}\left(U_{\alpha} \cap \Omega\right)$ for each $\alpha$; if $\Omega$ is not open, $C^{s}(\Omega, \mathbb{R})$ is the set of functions which have a $C^{s}$ extension to some neighborhood of $\Omega$.

A tangent vector $\xi$ at a point $a \in M$ is by definition a differential operator acting on functions, of the type

$$
C^{1}(\Omega, \mathbb{R}) \ni f \longmapsto \xi \cdot f=\sum_{1 \leqslant j \leqslant m} \xi_{j} \frac{\partial f}{\partial x_{j}}(a)
$$

in any local coordinate system $\left(x_{1}, \ldots, x_{m}\right)$ on an open set $\Omega \ni a$. We then simply write $\xi=\sum \xi_{j} \partial / \partial x_{j}$. For every $a \in \Omega$, the $n$-tuple $\left(\partial / \partial x_{j}\right)_{1 \leqslant j \leqslant m}$ is therefore a basis of the tangent space to $M$ at $a$, which we denote by $T_{M, a}$. The differential of a function $f$ at $a$ is the linear form on $T_{M, a}$ defined by

$$
d f_{a}(\xi)=\xi \cdot f=\sum \xi_{j} \partial f / \partial x_{j}(a), \quad \forall \xi \in T_{M, a}
$$

In particular $d x_{j}(\xi)=\xi_{j}$ and we may consequently write $d f=\sum\left(\partial f / \partial x_{j}\right) d x_{j}$. From this, we see that $\left(d x_{1}, \ldots, d x_{m}\right)$ is the dual basis of $\left(\partial / \partial x_{1}, \ldots, \partial / \partial x_{m}\right)$ in the cotangent space $T_{M, a}^{\star}$. The disjoint unions $T_{M}=\bigcup_{x \in M} T_{M, x}$ and $T_{M}^{\star}=\bigcup_{x \in M} T_{M, x}^{\star}$ are called the tangent and cotangent bundles of $M$.

If $\xi$ is a vector field of class $C^{s}$ over $\Omega$, that is, a map $x \mapsto \xi(x) \in T_{M, x}$ such that $\xi(x)=\sum \xi_{j}(x) \partial / \partial x_{j}$ has $C^{s}$ coefficients, and if $\eta$ is another vector field of class $C^{s}$ with $s \geqslant 1$, the Lie bracket $[\xi, \eta]$ is the vector field such that

$$
\begin{equation*}
[\xi, \eta] \cdot f=\xi \cdot(\eta \cdot f)-\eta \cdot(\xi \cdot f) \tag{1.2}
\end{equation*}
$$

In coordinates, it is easy to check that

$$
\begin{equation*}
[\xi, \eta]=\sum_{1 \leqslant j, k \leqslant m}\left(\xi_{j} \frac{\partial \eta_{k}}{\partial x_{j}}-\eta_{j} \frac{\partial \xi_{k}}{\partial x_{j}}\right) \frac{\partial}{\partial x_{k}} \tag{1.3}
\end{equation*}
$$

## § 1.B. Differential Forms

A differential form $u$ of degree $p$, or briefly a $p$-form over $M$, is a map $u$ on $M$ with values $u(x) \in \Lambda^{p} T_{M, x}^{\star}$. In a coordinate open set $\Omega \subset M$, a differential $p$-form can be written

$$
u(x)=\sum_{|I|=p} u_{I}(x) d x_{I},
$$

where $I=\left(i_{1}, \ldots, i_{p}\right)$ is a multi-index with integer components, $i_{1}<\ldots<i_{p}$ and $d x_{I}:=$ $d x_{i_{1}} \wedge \ldots \wedge d x_{i_{p}}$. The notation $|I|$ stands for the number of components of $I$, and is read length of $I$. For all integers $p=0,1, \ldots, m$ and $s \in \mathbb{N} \cup\{\infty\}, s \leqslant k$, we denote by $C^{s}\left(M, \Lambda^{p} T_{M}^{\star}\right)$ the space of differential $p$-forms of class $C^{s}$, i.e. with $C^{s}$ coefficients $u_{I}$. Several natural operations on differential forms can be defined.
$\S$ 1.B.1. Wedge Product. If $v(x)=\sum v_{J}(x) d x_{J}$ is a $q$-form, the wedge product of $u$ and $v$ is the form of degree $(p+q)$ defined by

$$
\begin{equation*}
u \wedge v(x)=\sum_{|I|=p,|J|=q} u_{I}(x) v_{J}(x) d x_{I} \wedge d x_{J} \tag{1.4}
\end{equation*}
$$

§ 1.B.2. Contraction by a tangent vector. A $p$-form $u$ can be viewed as an antisymmetric $p$-linear form on $T_{M}$. If $\xi=\sum \xi_{j} \partial / \partial x_{j}$ is a tangent vector, we define the contraction $\xi\lrcorner u$ to be the differential form of degree $p-1$ such that

$$
\begin{equation*}
(\xi\lrcorner u)\left(\eta_{1}, \ldots, \eta_{p-1}\right)=u\left(\xi, \eta_{1}, \ldots, \eta_{p-1}\right) \tag{1.5}
\end{equation*}
$$

for all tangent vectors $\eta_{j}$. Then $\left.(\xi, u) \longmapsto \xi\right\lrcorner u$ is bilinear and we find easily

$$
\left.\frac{\partial}{\partial x_{j}}\right\lrcorner d x_{I}= \begin{cases}0 & \text { if } j \notin I, \\ (-1)^{l-1} d x_{I \backslash\{j\}} & \text { if } j=i_{l} \in I\end{cases}
$$

A simple computation based on the above formula shows that contraction by a tangent vector is a derivation, i.e.

$$
\begin{equation*}
\left.\xi\lrcorner(u \wedge v)=(\xi\lrcorner u) \wedge v+(-1)^{\operatorname{deg} u} u \wedge(\xi\lrcorner v\right) . \tag{1.6}
\end{equation*}
$$

§ 1.B.3. Exterior derivative. This is the differential operator

$$
d: C^{s}\left(M, \Lambda^{p} T_{M}^{\star}\right) \longrightarrow C^{s-1}\left(M, \Lambda^{p+1} T_{M}^{\star}\right)
$$

defined in local coordinates by the formula

$$
\begin{equation*}
d u=\sum_{|I|=p, 1 \leqslant k \leqslant m} \frac{\partial u_{I}}{\partial x_{k}} d x_{k} \wedge d x_{I} . \tag{1.7}
\end{equation*}
$$

Alternatively, one can define $d u$ by its action on arbitrary vector fields $\xi_{0}, \ldots, \xi_{p}$ on $M$. The formula is as follows

$$
\begin{align*}
d u\left(\xi_{0}, \ldots, \xi_{p}\right) & =\sum_{0 \leqslant j \leqslant p}(-1)^{j} \xi_{j} \cdot u\left(\xi_{0}, \ldots, \widehat{\xi}_{j}, \ldots, \xi_{p}\right) \\
& +\sum_{0 \leqslant j<k \leqslant p}(-1)^{j+k} u\left(\left[\xi_{j}, \xi_{k}\right], \xi_{0}, \ldots, \widehat{\xi}_{j}, \ldots, \widehat{\xi}_{k}, \ldots, \xi_{p}\right)
\end{align*}
$$

The reader will easily check that (1.7) actually implies (1.7'). The advantage of (1.7') is that it does not depend on the choice of coordinates, thus $d u$ is intrinsically defined. The two basic properties of the exterior derivative (again left to the reader) are:

$$
\begin{align*}
& d(u \wedge v)=d u \wedge v+(-1)^{\operatorname{deg} u} u \wedge d v, \quad \text { (Leibnitz' rule ) }  \tag{1.8}\\
& d^{2}=0 .
\end{align*}
$$

A form $u$ is said to be closed if $d u=0$ and exact if $u$ can be written $u=d v$ for some form $v$.
§ 1.B.4. De Rham Cohomology Groups. Recall that a cohomological complex $K^{\bullet}=$ $\bigoplus_{p \in \mathbb{Z}}$ is a collection of modules $K^{p}$ over some ring, equipped with differentials, i.e., linear maps $d^{p}: K^{p} \rightarrow K^{p+1}$ such that $d^{p+1} \circ d^{p}=0$. The cocycle, coboundary and cohomology modules $Z^{p}\left(K^{\bullet}\right), B^{p}\left(K^{\bullet}\right)$ and $H^{p}\left(K^{\bullet}\right)$ are defined respectively by

$$
\begin{cases}Z^{p}\left(K^{\bullet}\right)=\operatorname{Ker} d^{p}: K^{p} \rightarrow K^{p+1}, & Z^{p}\left(K^{\bullet}\right) \subset K^{p},  \tag{1.10}\\ B^{p}\left(K^{\bullet}\right)=\operatorname{Im} d^{p-1}: K^{p-1} \rightarrow K^{p}, & B^{p}\left(K^{\bullet}\right) \subset Z^{p}\left(K^{\bullet}\right) \subset K^{p}, \\ H^{p}\left(K^{\bullet}\right)=Z^{p}\left(K^{\bullet}\right) / B^{p}\left(K^{\bullet}\right) . & \end{cases}
$$

Now, let $M$ be a differentiable manifold, say of class $\mathscr{C}^{\infty}$ for simplicity. The De Rham complex of $M$ is defined to be the complex $K^{p}=\mathscr{C}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star}\right)$ of smooth differential forms, together with the exterior derivative $d^{p}=d$ as differential, and $K^{p}=\{0\}, d^{p}=0$ for $p<0$. We denote by $Z^{p}(M, \mathbb{R})$ the cocycles (closed $p$-forms) and by $B^{p}(M, \mathbb{R})$ the coboundaries (exact $p$-forms). By convention $B^{0}(M, \mathbb{R})=\{0\}$. The De Rham cohomology group of $M$ in degree $p$ is

$$
\begin{equation*}
H_{\mathrm{DR}}^{p}(M, \mathbb{R})=Z^{p}(M, \mathbb{R}) / B^{p}(M, \mathbb{R}) \tag{1.11}
\end{equation*}
$$

When no confusion with other types of cohomology groups may occur, we sometimes denote these groups simply by $H^{p}(M, \mathbb{R})$. The symbol $\mathbb{R}$ is used here to stress that we are considering real valued $p$-forms; of course one can introduce a similar group $H_{\mathrm{DR}}^{p}(M, \mathbb{C})$ for complex valued forms, i.e. forms with values in $\mathbb{C} \otimes \Lambda^{p} T_{M}^{\star}$. Then $H_{\mathrm{DR}}^{p}(M, \mathbb{C})=$ $\mathbb{C} \otimes H_{\mathrm{DR}}^{p}(M, \mathbb{R})$ is the complexification of the real De Rham cohomology group. It is clear that $H_{\mathrm{DR}}^{0}(M, \mathbb{R})$ can be identified with the space of locally constant functions on $M$, thus

$$
H_{\mathrm{DR}}^{0}(M, \mathbb{R})=\mathbb{R}^{\pi_{0}(X)}
$$

where $\pi_{0}(X)$ denotes the set of connected components of $M$.
Similarly, we introduce the De Rham cohomology groups with compact support

$$
\begin{equation*}
H_{\mathrm{DR}, c}^{p}(M, \mathbb{R})=Z_{c}^{p}(M, \mathbb{R}) / B_{c}^{p}(M, \mathbb{R}) \tag{1.12}
\end{equation*}
$$

associated with the De Rham complex $K^{p}=\mathscr{C}_{c}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star}\right)$ of smooth differential forms with compact support.
§ 1.B.5. Pull-Back. If $F: M \longrightarrow M^{\prime}$ is a differentiable map to another manifold $M^{\prime}$, $\operatorname{dim}_{\mathbb{R}} M^{\prime}=m^{\prime}$, and if $v(y)=\sum v_{J}(y) d y_{J}$ is a differential $p$-form on $M^{\prime}$, the pull-back $F^{\star} v$ is the differential $p$-form on $M$ obtained after making the substitution $y=F(x)$ in $v$, i.e.

$$
\begin{equation*}
F^{\star} v(x)=\sum v_{I}(F(x)) d F_{i_{1}} \wedge \ldots \wedge d F_{i_{p}} \tag{1.13}
\end{equation*}
$$

If we have a second map $G: M^{\prime} \longrightarrow M^{\prime \prime}$ and if $w$ is a differential form on $M^{\prime \prime}$, then $F^{\star}\left(G^{\star} w\right)$ is obtained by means of the substitutions $z=G(y), y=F(x)$, thus

$$
\begin{equation*}
F^{\star}\left(G^{\star} w\right)=(G \circ F)^{\star} w \tag{1.14}
\end{equation*}
$$

Moreover, we always have $d\left(F^{\star} v\right)=F^{\star}(d v)$. It follows that the pull-back $F^{\star}$ is closed if $v$ is closed and exact if $v$ is exact. Therefore $F^{\star}$ induces a morphism on the quotient spaces

$$
\begin{equation*}
F^{\star}: H_{\mathrm{DR}}^{p}\left(M^{\prime}, \mathbb{R}\right) \longrightarrow H_{\mathrm{DR}}^{p}(M, \mathbb{R}) \tag{1.15}
\end{equation*}
$$

## § 1.C. Integration of Differential Forms

A manifold $M$ is orientable if and only if there exists an atlas $\left(\tau_{\alpha}\right)$ such that all transition maps $\tau_{\alpha \beta}$ preserve the orientation, i.e. have positive jacobian determinants. Suppose that $M$ is oriented, that is, equipped with such an atlas. If $u(x)=f\left(x_{1}, \ldots, x_{m}\right) d x_{1} \wedge$ $\ldots \wedge d x_{m}$ is a continuous form of maximum degree $m=\operatorname{dim}_{\mathbb{R}} M$, with compact support in a coordinate open set $\Omega$, we set

$$
\begin{equation*}
\int_{M} u=\int_{\mathbb{R}^{m}} f\left(x_{1}, \ldots, x_{m}\right) d x_{1} \ldots d x_{m} \tag{1.16}
\end{equation*}
$$

By the change of variable formula, the result is independent of the choice of coordinates, provided we consider only coordinates corresponding to the given orientation. When $u$ is an arbitrary form with compact support, the definition of $\int_{M} u$ is easily extended by means of a partition of unity with respect to coordinate open sets covering Supp $u$. Let $F: M \longrightarrow M^{\prime}$ be a diffeomorphism between oriented manifolds and $v$ a volume form on $M^{\prime}$. The change of variable formula yields

$$
\begin{equation*}
\int_{M} F^{\star} v= \pm \int_{M^{\prime}} v \tag{1.17}
\end{equation*}
$$

according whether $F$ preserves orientation or not.
We now state Stokes' formula, which is basic in many contexts. Let $K$ be a compact subset of $M$ with piecewise $C^{1}$ boundary. By this, we mean that for each point $a \in \partial K$ there are coordinates $\left(x_{1}, \ldots, x_{m}\right)$ on a neighborhood $V$ of $a$, centered at $a$, such that

$$
K \cap V=\left\{x \in V ; x_{1} \leqslant 0, \ldots, x_{l} \leqslant 0\right\}
$$

for some index $l \geqslant 1$. Then $\partial K \cap V$ is a union of smooth hypersurfaces with piecewise $C^{1}$ boundaries:

$$
\partial K \cap V=\bigcup_{1 \leqslant j \leqslant l}\left\{x \in V ; x_{1} \leqslant 0, \ldots, x_{j}=0, \ldots, x_{l} \leqslant 0\right\}
$$

At points of $\partial K$ where $x_{j}=0$, then $\left(x_{1}, \ldots, \widehat{x_{j}}, \ldots, x_{m}\right)$ define coordinates on $\partial K$. We take the orientation of $\partial K$ given by these coordinates or the opposite one, according to the sign $(-1)^{j-1}$. For any differential form $u$ of class $C^{1}$ and degree $m-1$ on $M$, we then have
(1.18) Stokes' formula. $\int_{\partial K} u=\int_{K} d u$.

The formula is easily checked by an explicit computation when $u$ has compact support in $V$ : indeed if $u=\sum_{1 \leqslant j \leqslant n} u_{j} d x_{1} \wedge \ldots \widehat{d x_{j}} \ldots d x_{m}$ and $\partial_{j} K \cap V$ is the part of $\partial K \cap V$ where $x_{j}=0$, a partial integration with respect to $x_{j}$ yields

$$
\begin{aligned}
& \int_{\partial_{j} K \cap V} u_{j} d x_{1} \wedge \ldots \widehat{d x_{j}} \ldots d x_{m}=\int_{V} \frac{\partial u_{j}}{\partial x_{j}} d x_{1} \wedge \ldots d x_{m}, \\
& \int_{\partial K \cap V} u=\sum_{1 \leqslant j \leqslant m}(-1)^{j-1} \int_{\partial_{j} K \cap V} u_{j} d x_{1} \wedge \ldots \widehat{d x_{j}} \ldots \wedge d x_{m}=\int_{V} d u .
\end{aligned}
$$

The general case follows by a partition of unity. In particular, if $u$ has compact support in $M$, we find $\int_{M} d u=0$ by choosing $K \supset \operatorname{Supp} u$.

## §1.D. Homotopy Formula and Poincaré Lemma

Let $u$ be a differential form on $[0,1] \times M$. For $(t, x) \in[0,1] \times M$, we write

$$
u(t, x)=\sum_{|I|=p} u_{I}(t, x) d x_{I}+\sum_{|J|=p-1} \widetilde{u}_{J}(t, x) d t \wedge d x_{J} .
$$

We define an operator

$$
\begin{align*}
& K: C^{s}\left([0,1] \times M, \Lambda^{p} T_{[0,1] \times M}^{\star}\right) \longrightarrow C^{s}\left(M, \Lambda^{p-1} T_{M}^{\star}\right) \\
& K u(x)=\sum_{|J|=p-1}\left(\int_{0}^{1} \widetilde{u}_{J}(t, x) d t\right) d x_{J} \tag{1.19}
\end{align*}
$$

and say that $K u$ is the form obtained by integrating $u$ along $[0,1]$. A computation of the operator $d K+K d$ shows that all terms involving partial derivatives $\partial \widetilde{u}_{J} / \partial x_{k}$ cancel, hence

$$
K d u+d K u=\sum_{|I|=p}\left(\int_{0}^{1} \frac{\partial u_{I}}{\partial t}(t, x) d t\right) d x_{I}=\sum_{|I|=p}\left(u_{I}(1, x)-u_{I}(0, x)\right) d x_{I}
$$

$$
\begin{equation*}
K d u+d K u=i_{1}^{\star} u-i_{0}^{\star} u \tag{1.20}
\end{equation*}
$$

where $i_{t}: M \rightarrow[0,1] \times M$ is the injection $x \mapsto(t, x)$.
(1.20) Corollary. Let $F, G: M \longrightarrow M^{\prime}$ be $\mathscr{C}^{\infty}$ maps. Suppose that $F, G$ are smoothly homotopic, i.e. that there exists a $\mathscr{C}^{\infty}$ map $H:[0,1] \times M \longrightarrow M^{\prime}$ such that $H(0, x)=$ $F(x)$ and $H(1, x)=G(x)$. Then

$$
F^{\star}=G^{\star}: H_{\mathrm{DR}}^{p}\left(M^{\prime}, \mathbb{R}\right) \longrightarrow H_{\mathrm{DR}}^{p}(M, \mathbb{R}) .
$$

Proof. If $v$ is a $p$-form on $M^{\prime}$, then

$$
\begin{aligned}
G^{\star} v-F^{\star} v & =\left(H \circ i_{1}\right)^{\star} v-\left(H \circ i_{0}\right)^{\star} v=i_{1}^{\star}\left(H^{\star} v\right)-i_{0}^{\star}\left(H^{\star} v\right) \\
& =d\left(K H^{\star} v\right)+K H^{\star}(d v)
\end{aligned}
$$

by (1.20) applied to $u=H^{\star} v$. If $v$ is closed, then $F^{\star} v$ and $G^{\star} v$ differ by an exact form, so they define the same class in $H_{\mathrm{DR}}^{p}(M, \mathbb{R})$.
(1.21) Corollary. If the manifold $M$ is contractible, i.e. if there is a smooth homotopy $H:[0,1] \times M \rightarrow M$ from a constant map $F: M \rightarrow\left\{x_{0}\right\}$ to $G=\operatorname{Id}_{X}$, then $H_{\mathrm{DR}}^{0}(M, \mathbb{R})=$ $\mathbb{R}$ and $H_{\mathrm{DR}}^{p}(M, \mathbb{R})=0$ for $p \geqslant 1$.

Proof. $F^{\star}$ is clearly zero in degree $p \geqslant 1$, while $F^{\star}: H_{\mathrm{DR}}^{0}(M, \mathbb{R}) \xrightarrow{\simeq} \mathbb{R}$ is induced by the evaluation map $u \mapsto u\left(x_{0}\right)$. The conclusion then follows from the equality $F^{\star}=G^{\star}=\mathrm{Id}$ on cohomology groups.
(1.22) Poincaré lemma. Let $\Omega \subset \mathbb{R}^{m}$ be a starshaped open set. If a form $v=$ $\sum v_{I} d x_{I} \in C^{s}\left(\Omega, \Lambda^{p} T_{\Omega}^{\star}\right), p \geqslant 1$, satisfies $d v=0$, there exists a form $u \in C^{s}\left(\Omega, \Lambda^{p-1} T_{\Omega}^{\star}\right)$ such that $d u=v$.

Proof. Let $H(t, x)=t x$ be the homotopy between the identity map $\Omega \rightarrow \Omega$ and the constant map $\Omega \rightarrow\{0\}$. By the above formula

$$
d\left(K H^{\star} v\right)=G^{\star} v-F^{\star} v= \begin{cases}v-v(0) & \text { if } p=0 \\ v & \text { if } p \geqslant 1\end{cases}
$$

Hence $u=K H^{\star} v$ is the ( $p-1$ )-form we are looking for. An explicit computation based on (1.19) easily gives

$$
\begin{equation*}
u(x)=\sum_{\substack{|I|=p \\ 1 \leqslant k \leqslant p}}\left(\int_{0}^{1} t^{p-1} v_{I}(t x) d t\right)(-1)^{k-1} x_{i_{k}} d x_{i_{1}} \wedge \ldots \widehat{d x_{i_{k}}} \ldots \wedge d x_{i_{p}} \tag{1.23}
\end{equation*}
$$

## § 2. Currents on Differentiable Manifolds

## $\S$ 2.A. Definition and Examples

Let $M$ be a $\mathscr{C}^{\infty}$ differentiable manifold, $m=\operatorname{dim}_{\mathbb{R}} M$. All the manifolds considered in Sect. 2 will be assumed to be oriented. We first introduce a topology on the space of differential forms $C^{s}\left(M, \Lambda^{p} T_{M}^{\star}\right)$. Let $\Omega \subset M$ be a coordinate open set and $u$ a $p$-form on $M$, written $u(x)=\sum u_{I}(x) d x_{I}$ on $\Omega$. To every compact subset $L \subset \Omega$ and every integer $s \in \mathbb{N}$, we associate a seminorm

$$
\begin{equation*}
p_{L}^{s}(u)=\sup _{x \in L} \max _{|I|=p,|\alpha| \leqslant s}\left|D^{\alpha} u_{I}(x)\right|, \tag{2.1}
\end{equation*}
$$

where $\alpha=\left(\alpha_{1}, \ldots, \alpha_{m}\right)$ runs over $\mathbb{N}^{m}$ and $D^{\alpha}=\partial^{|\alpha|} / \partial x_{1}^{\alpha_{1}} \ldots \partial x_{m}^{\alpha_{m}}$ is a derivation of order $|\alpha|=\alpha_{1}+\cdots+\alpha_{m}$. This type of multi-index, which will always be denoted by Greek letters, should not be confused with multi-indices of the type $I=\left(i_{1}, \ldots, i_{p}\right)$ introduced in Sect. 1.
(2.2) Definition. We introduce as follows spaces of $p$-forms on manifolds.
a) We denote by $\mathscr{C}^{p}(M)\left(\right.$ resp. $\left.{ }^{s} \mathscr{C}^{p}(M)\right)$ the space $\mathscr{C}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star}\right)\left(\right.$ resp. the space $\left.C^{s}\left(M, \Lambda^{p} T_{M}^{\star}\right)\right)$, equipped with the topology defined by all seminorms $p_{L}^{s}$ when $s, L, \Omega$ vary (resp. when $L, \Omega$ vary $)$.
b) If $K \subset M$ is a compact subset, $\mathscr{D}^{p}(K)$ will denote the subspace of elements $u \in \mathscr{E}^{p}(M)$ with support contained in $K$, together with the induced topology; $\mathscr{D}^{p}(M)$ will stand for the set of all elements with compact support, i.e. $\mathscr{D}^{p}(M):=\bigcup_{K} \mathscr{D}^{p}(K)$.
c) The spaces of $C^{s}$-forms $s^{s} \mathscr{D}^{p}(K)$ and ${ }^{s} \mathscr{D}^{p}(M)$ are defined similarly.

Since our manifolds are assumed to be separable, the topology of $\mathscr{E}^{p}(M)$ can be defined by means of a countable set of seminorms $p_{L}^{s}$, hence $\mathscr{E}^{p}(M)$ (and likewise ${ }^{s} \mathscr{C}^{p}(M)$ ) is a Fréchet space. The topology of ${ }^{s} \mathscr{D}^{p}(K)$ is induced by any finite set of seminorms $p_{K_{j}}^{s}$ such that the compact sets $K_{j}$ cover $K$; hence ${ }^{s} \mathscr{D}^{p}(K)$ is a Banach space. It should be observed however that $\mathscr{D}^{p}(M)$ is not a Fréchet space; in fact $\mathscr{D}^{p}(M)$ is dense in $\mathscr{E}^{p}(M)$ and thus non complete for the induced topology. According to [De Rham 1955] spaces of currents are defined as the topological duals of the above spaces, in analogy with the usual definition of distributions.
(2.3) Definition. The space of currents of dimension $p$ (or degree $m-p$ ) on $M$ is the space $\mathscr{D}_{p}^{\prime}(M)$ of linear forms $T$ on $\mathscr{D}^{p}(M)$ such that the restriction of $T$ to all subspaces $\mathscr{D}^{p}(K), K \subset \subset M$, is continuous. The degree is indicated by raising the index, hence we set

$$
\mathscr{D}^{\prime m-p}(M)=\mathscr{D}_{p}^{\prime}(M):=\text { topological dual }\left(\mathscr{D}^{p}(M)\right)^{\prime} .
$$

The space ${ }^{s} \mathscr{D}_{p}^{\prime}(M)={ }^{s} \mathscr{D}^{\prime m-p}(M):=\left({ }^{s} \mathscr{D}^{p}(M)\right)^{\prime}$ is defined similarly and is called the space of currents of order s on $M$.

In the sequel, we let $\langle T, u\rangle$ be the pairing between a current $T$ and a test form $u \in \mathscr{D}^{p}(M)$. It is clear that ${ }^{s} \mathscr{D}_{p}^{\prime}(M)$ can be identified with the subspace of currents $T \in \mathscr{D}_{p}^{\prime}(M)$ which are continuous for the seminorm $p_{K}^{s}$ on $\mathscr{D}^{p}(K)$ for every compact set $K$ contained in a coordinate patch $\Omega$. The support of $T$, denoted Supp $T$, is the smallest closed subset $A \subset M$ such that the restriction of $T$ to $\mathscr{D}^{p}(M \backslash A)$ is zero. The topological dual $\mathscr{E}_{p}^{\prime}(M)$ can be identified with the set of currents of $\mathscr{D}_{p}^{\prime}(M)$ with compact support: indeed, let $T$ be a linear form on $\mathscr{E}^{p}(M)$ such that

$$
|\langle T, u\rangle| \leqslant C \max \left\{p_{K_{j}}^{s}(u)\right\}
$$

for some $s \in \mathbb{N}, C \geqslant 0$ and a finite number of compact sets $K_{j}$; it follows that $\operatorname{Supp} T \subset$ $\bigcup K_{j}$. Conversely let $T \in \mathscr{D}_{p}^{\prime}(M)$ with support in a compact set $K$. Let $K_{j}$ be compact patches such that $K$ is contained in the interior of $\bigcup K_{j}$ and $\psi \in \mathscr{D}(M)$ equal to 1 on $K$ with $\operatorname{Supp} \psi \subset \bigcup K_{j}$. For $u \in \mathscr{E}^{p}(M)$, we define $\langle T, u\rangle=\langle T, \psi u\rangle$; this is independent of $\psi$ and the resulting $T$ is clearly continuous on $\mathscr{C}^{p}(M)$. The terminology used for the dimension and degree of a current is justified by the following two examples.
(2.4) Example. Let $Z \subset M$ be a closed oriented submanifold of $M$ of dimension $p$ and class $C^{1} ; Z$ may have a boundary $\partial Z$. The current of integration over $Z$, denoted $[Z]$, is defined by

$$
\langle[Z], u\rangle=\int_{Z} u, \quad u \in{ }^{0} \mathscr{D}^{p}(M)
$$

It is clear that $[Z]$ is a current of order 0 on $M$ and that $\operatorname{Supp}[Z]=Z$. Its dimension is $p=\operatorname{dim} Z$.
(2.5) Example. If $f$ is a differential form of degree $q$ on $M$ with $L_{\text {loc }}^{1}$ coefficients, we can associate to $f$ the current of dimension $m-q$ :

$$
\left\langle T_{f}, u\right\rangle=\int_{M} f \wedge u, \quad u \in{ }^{0} \mathscr{D}^{m-q}(M)
$$

$T_{f}$ is of degree $q$ and of order 0 . The correspondence $f \longmapsto T_{f}$ is injective. In the same way $L_{\text {loc }}^{1}$ functions on $\mathbb{R}^{m}$ are identified to distributions, we will identify $f$ with its image $T_{f} \in{ }^{0} \mathscr{D}^{\prime} q(M)={ }^{0} \mathscr{D}_{m-q}^{\prime}(M)$.

## §2.B. Exterior Derivative and Wedge Product

$\S$ 2.B.1. Exterior Derivative. Many of the operations available for differential forms can be extended to currents by simple duality arguments. Let $T \in{ }^{s} \mathscr{D}^{\prime q}(M)={ }^{s} \mathscr{D}_{m-p}^{\prime}(M)$. The exterior derivative

$$
d T \in{ }^{s+1} \mathscr{D}^{\prime q+1}(M)={ }^{s+1} \mathscr{D}_{m-q-1}^{\prime}
$$

is defined by

$$
\begin{equation*}
\langle d T, u\rangle=(-1)^{q+1}\langle T, d u\rangle, \quad u \in^{s+1} \mathscr{D}^{m-q-1}(M) \tag{2.6}
\end{equation*}
$$

The continuity of the linear form $d T$ on ${ }^{s+1} \mathscr{D}^{m-q-1}(M)$ follows from the continuity of the map $d:{ }^{s+1} \mathscr{D}^{m-q-1}(K) \longrightarrow{ }^{s} \mathscr{D}^{m-q}(K)$. For all forms $f \in{ }^{1} \mathscr{C}^{q}(M)$ and $u \in \mathscr{D}^{m-q-1}(M)$, Stokes' formula implies

$$
0=\int_{M} d(f \wedge u)=\int_{M} d f \wedge u+(-1)^{q} f \wedge d u
$$

thus in example (2.5) one actually has $d T_{f}=T_{d f}$ as it should be. In example (2.4), another application of Stokes' formula yields $\int_{Z} d u=\int_{\partial Z} u$, therefore $\langle[Z], d u\rangle=\langle[\partial Z], u\rangle$ and

$$
\begin{equation*}
d[Z]=(-1)^{m-p+1}[\partial Z] . \tag{2.7}
\end{equation*}
$$

$\S$ 2.B.2. Wedge Product. For $T \in{ }^{s} \mathscr{D}^{\prime q}(M)$ and $g \in{ }^{s} \mathscr{E}^{r}(M)$, the wedge product $T \wedge g \in{ }^{s} \mathscr{D}^{\prime q+r}(M)$ is defined by

$$
\begin{equation*}
\langle T \wedge g, u\rangle=\langle T, g \wedge u\rangle, \quad u \in{ }^{s} \mathscr{D}^{m-q-r}(M) \tag{2.8}
\end{equation*}
$$

This definition is licit because $u \mapsto g \wedge u$ is continuous in the $C^{s}$-topology. The relation

$$
d(T \wedge g)=d T \wedge g+(-1)^{\operatorname{deg} T} T \wedge d g
$$

is easily verified from the definitions.
(2.9) Proposition. Let $\left(x_{1}, \ldots, x_{m}\right)$ be a coordinate system on an open subset $\Omega \subset M$. Every current $T \in{ }^{s} \mathscr{D}^{\prime}(M)$ of degree $q$ can be written in a unique way

$$
T=\sum_{|I|=q} T_{I} d x_{I} \quad \text { on } \quad \Omega,
$$

where $T_{I}$ are distributions of order $s$ on $\Omega$, considered as currents of degree 0 .
Proof. If the result is true, for all $f \in{ }^{s} \mathscr{D}^{0}(\Omega)$ we must have

$$
\left\langle T, f d x_{\mathrm{C}_{I}}\right\rangle=\left\langle T_{I}, d x_{I} \wedge f d x_{\mathrm{C}_{I}}\right\rangle=\varepsilon(I, \mathrm{CI})\left\langle T_{I}, f d x_{1} \wedge \ldots \wedge d x_{m}\right\rangle
$$

where $\varepsilon(I, \complement I)$ is the signature of the permutation $(1, \ldots, m) \longmapsto(I, C I)$. Conversely, this can be taken as a definition of the coefficient $T_{I}$ :

$$
\begin{equation*}
T_{I}(f)=\left\langle T_{I}, f d x_{1} \wedge \ldots \wedge d x_{m}\right\rangle:=\varepsilon(I, \complement I)\left\langle T, f d x_{C_{I}}\right\rangle, \quad f \in{ }^{s} \mathscr{D}^{0}(\Omega) . \tag{2.10}
\end{equation*}
$$

Then $T_{I}$ is a distribution of order $s$ and it is easy to check that $T=\sum T_{I} d x_{I}$.
In particular, currents of order 0 on $M$ can be considered as differential forms with measure coefficients. In order to unify the notations concerning forms and currents, we set

$$
\langle T, u\rangle=\int_{M} T \wedge u
$$

whenever $T \in{ }^{s} \mathscr{D}_{p}^{\prime}(M)={ }^{s} \mathscr{D}^{\prime m-p}(M)$ and $u \in{ }^{s} \mathscr{C}^{p}(M)$ are such that $\operatorname{Supp} T \cap \operatorname{Supp} u$ is compact. This convention is made so that the notation becomes compatible with the identification of a form $f$ to the current $T_{f}$.

## § 2.C. Direct and Inverse Images

$\S$ 2.C.1. Direct Images. Assume now that $M_{1}, M_{2}$ are oriented differentiable manifolds of respective dimensions $m_{1}, m_{2}$, and that

$$
\begin{equation*}
F: M_{1} \longrightarrow M_{2} \tag{2.11}
\end{equation*}
$$

is a $\mathscr{C}^{\infty}$ map. The pull-back morphism

$$
\begin{equation*}
{ }^{s} \mathscr{D}^{p}\left(M_{2}\right) \longrightarrow{ }^{s} \mathscr{E}^{p}\left(M_{1}\right), \quad u \longmapsto F^{\star} u \tag{2.12}
\end{equation*}
$$

is continuous in the $C^{s}$ topology and we have $\operatorname{Supp} F^{\star} u \subset F^{-1}(\operatorname{Supp} u)$, but in general Supp $F^{\star} u$ is not compact. If $T \in{ }^{s} \mathscr{D}_{p}^{\prime}\left(M_{1}\right)$ is such that the restriction of $F$ to $\operatorname{Supp} T$ is proper, i.e. if $\operatorname{Supp} T \cap F^{-1}(K)$ is compact for every compact subset $K \subset M_{2}$, then the linear form $u \longmapsto\left\langle T, F^{\star} u\right\rangle$ is well defined and continuous on ${ }^{s} \mathscr{D}^{p}\left(M_{2}\right)$. There exists therefore a unique current denoted $F_{\star} T \in{ }^{s} \mathscr{D}_{p}^{\prime}\left(M_{2}\right)$, called the direct image of $T$ by $F$, such that

$$
\begin{equation*}
\left\langle F_{\star} T, u\right\rangle=\left\langle T, F^{\star} u\right\rangle, \quad \forall u \in{ }^{s} \mathscr{D}^{p}\left(M_{2}\right) . \tag{2.13}
\end{equation*}
$$

We leave the straightforward proof of the following properties to the reader.
(2.14) Theorem. For every $T \in{ }^{s} \mathscr{D}_{p}^{\prime}\left(M_{1}\right)$ such that $F_{\lceil\text {Supp } T}$ is proper, the direct image $F_{\star} T \in{ }^{s} \mathscr{D}_{p}^{\prime}\left(M_{2}\right)$ is such that
a) $\operatorname{Supp} F_{\star} T \subset F(\operatorname{Supp} T)$;
b) $d\left(F_{\star} T\right)=F_{\star}(d T)$;
c) $F_{\star}\left(T \wedge F^{\star} g\right)=\left(F_{\star} T\right) \wedge g, \quad \forall g \in^{s} \mathscr{E}^{q}\left(M_{2}, \mathbb{R}\right)$;
d) If $G: M_{2} \longrightarrow M_{3}$ is a $\mathscr{C}^{\infty}$ map such that $(G \circ F)_{\mid \operatorname{Supp} T}$ is proper, then

$$
G_{\star}\left(F_{\star} T\right)=(G \circ F)_{\star} T .
$$

(2.15) Special case. Assume that $F$ is a submersion, i.e. that $F$ is surjective and that for every $x \in M_{1}$ the differential map $d_{x} F: T_{M_{1}, x} \longrightarrow T_{M_{2}, F(x)}$ is surjective. Let $g$ be a differential form of degree $q$ on $M_{1}$, with $L_{\text {loc }}^{1}$ coefficients, such that $F_{\text {โSupp } g}$ is proper. We claim that $F_{\star} g \in{ }^{0} \mathscr{D}_{m_{1}-q}^{\prime}\left(M_{2}\right)$ is the form of degree $q-\left(m_{1}-m_{2}\right)$ obtained from $g$ by integration along the fibers of $F$, also denoted

$$
F_{\star} g(y)=\int_{z \in F^{-1}(y)} g(z)
$$



Fig. I-2 Local description of a submersion as a projection.

In fact, this assertion is equivalent to the following generalized form of Fubini's theorem:

$$
\int_{M_{1}} g \wedge F^{\star} u=\int_{y \in M_{2}}\left(\int_{z \in F^{-1}(y)} g(z)\right) \wedge u(y), \quad \forall u \in{ }^{0} \mathscr{D}^{m_{1}-q}\left(M_{2}\right)
$$

By using a partition of unity on $M_{1}$ and the constant rank theorem, the verification of this formula is easily reduced to the case where $M_{1}=A \times M_{2}$ and $F=\mathrm{pr}_{2}$, cf. Fig. 2. The fibers $F^{-1}(y) \simeq A$ have to be oriented in such a way that the orientation of $M_{1}$ is the product of the orientation of $A$ and $M_{2}$. Let us write $r=\operatorname{dim} A=m_{1}-m_{2}$ and let $z=(x, y) \in A \times M_{2}$ be any point of $M_{1}$. The above formula becomes

$$
\int_{A \times M_{2}} g(x, y) \wedge u(y)=\int_{y \in M_{2}}\left(\int_{x \in A} g(x, y)\right) \wedge u(y)
$$

where the direct image of $g$ is computed from $g=\sum g_{I, J}(x, y) d x_{I} \wedge d y_{J},|I|+|J|=q$, by the formula

$$
\begin{align*}
F_{\star} g(y) & =\int_{x \in A} g(x, y)  \tag{2.16}\\
& =\sum_{|J|=q-r}\left(\int_{x \in A} g_{(1, \ldots, r), J}(x, y) d x_{1} \wedge \ldots \wedge d x_{r}\right) d y_{J} .
\end{align*}
$$

In this situation, we see that $F_{\star} g$ has $L_{\mathrm{loc}}^{1}$ coefficients on $M_{2}$ if $g$ is $L_{\mathrm{loc}}^{1}$ on $M_{1}$, and that the map $g \longmapsto F_{\star} g$ is continuous in the $C^{s}$ topology.
(2.17) Remark. If $F: M_{1} \longrightarrow M_{2}$ is a diffeomorphism, then we have $F_{\star} g= \pm\left(F^{-1}\right)^{\star} g$ according whether $F$ preserves the orientation or not. In fact formula (1.17) gives

$$
\left\langle F_{\star} g, u\right\rangle=\int_{M_{1}} g \wedge F^{\star} u= \pm \int_{M_{2}}\left(F^{-1}\right)^{\star}\left(g \wedge F^{\star} u\right)= \pm \int_{M_{2}}\left(F^{-1}\right)^{\star} g \wedge u
$$

$\S$ 2.C.2. Inverse Images. Assume that $F: M_{1} \longrightarrow M_{2}$ is a submersion. As a consequence of the continuity statement after (2.16), one can always define the inverse image $F^{\star} T \in$ ${ }^{s} \mathscr{D}^{\prime q}\left(M_{1}\right)$ of a current $T \in{ }^{s} \mathscr{D}^{\prime q}\left(M_{2}\right)$ by

$$
\left\langle F^{\star} T, u\right\rangle=\left\langle T, F_{\star} u\right\rangle, \quad u \in{ }^{s} \mathscr{D}^{q+m_{1}-m_{2}}\left(M_{1}\right) .
$$

Then $\operatorname{dim} F^{\star} T=\operatorname{dim} T+m_{1}-m_{2}$ and Th. 2.14 yields the formulas:

$$
\begin{equation*}
d\left(F^{\star} T\right)=F^{\star}(d T), \quad F^{\star}(T \wedge g)=F^{\star} T \wedge F^{\star} g, \quad \forall g \in^{s} \mathscr{D} \bullet\left(M_{2}\right) \tag{2.18}
\end{equation*}
$$

Take in particular $T=[Z]$, where $Z$ is an oriented $C^{1}$-submanifold of $M_{2}$. Then $F^{-1}(Z)$ is a submanifold of $M_{1}$ and has a natural orientation given by the isomorphism

$$
T_{M_{1}, x} / T_{F^{-1}(Z), x} \longrightarrow T_{M_{2}, F(x)} / T_{Z, F(x)},
$$

induced by $d_{x} F$ at every point $x \in Z$. We claim that

$$
\begin{equation*}
F^{\star}[Z]=\left[F^{-1}(Z)\right] . \tag{2.19}
\end{equation*}
$$

Indeed, we have to check that $\int_{Z} F_{\star} u=\int_{F^{-1}(Z)} u$ for every $u \in{ }^{s} \mathscr{D}^{\bullet}\left(M_{1}\right)$. By using a partition of unity on $M_{1}$, we may again assume $M_{1}=A \times M_{2}$ and $F=\operatorname{pr}_{2}$. The above equality can be written

$$
\int_{y \in Z} F_{\star} u(y)=\int_{(x, y) \in A \times Z} u(x, y) .
$$

This follows precisely from (2.16) and Fubini's theorem.
§ 2.C.3. Weak Topology. The weak topology on $\mathscr{D}_{p}^{\prime}(M)$ is the topology defined by the collection of seminorms $T \longmapsto|\langle T, f\rangle|$ for all $f \in \mathscr{D}^{p}(M)$. With respect to the weak topology, all the operations

$$
\begin{equation*}
T \longmapsto d T, \quad T \longmapsto T \wedge g, \quad T \longmapsto F_{\star} T, \quad T \longmapsto F^{\star} T \tag{2.20}
\end{equation*}
$$

defined above are continuous. A set $B \subset \mathscr{D}_{p}^{\prime}(M)$ is bounded for the weak topology (weakly bounded for short) if and only if $\langle T, f\rangle$ is bounded when $T$ runs over $B$, for every fixed $f \in \mathscr{D}^{p}(M)$. The standard Banach-Alaoglu theorem implies that every weakly bounded closed subset $B \subset \mathscr{D}_{p}^{\prime}(M)$ is weakly compact.

## $\S$ 2.D. Tensor Products, Homotopies and Poincaré Lemma

$\S$ 2.D.1. Tensor Products. If $S, T$ are currents on manifolds $M, M^{\prime}$ there exists a unique current on $M \times M^{\prime}$, denoted $S \otimes T$ and defined in a way analogous to the tensor product of distributions, such that for all $u \in \mathscr{D}^{\bullet}(M)$ and $v \in \mathscr{D}^{\bullet}\left(M^{\prime}\right)$

$$
\begin{equation*}
\left\langle S \otimes T, \operatorname{pr}_{1}^{\star} u \wedge \operatorname{pr}_{2}^{\star} v\right\rangle=(-1)^{\operatorname{deg} T \operatorname{deg} u}\langle S, u\rangle\langle T, v\rangle . \tag{2.21}
\end{equation*}
$$

One verifies easily that $d(S \otimes T)=d S \otimes T+(-1)^{\operatorname{deg} S} S \otimes d T$.
§ 2.D.2. Homotopy Formula. Assume that $H:[0,1] \times M_{1} \longrightarrow M_{2}$ is a $\mathscr{C}^{\infty}$ homotopy from $F(x)=H(0, x)$ to $G(x)=H(1, x)$ and that $T \in \mathscr{D}_{\bullet}^{\prime}\left(M_{1}\right)$ is a current such that $H_{\lceil[0,1] \times \operatorname{Supp} T}$ is proper. If $[0,1]$ is considered as the current of degree 0 on $\mathbb{R}$ associated to its characteristic function, we find $d[0,1]=\delta_{0}-\delta_{1}$, thus

$$
\begin{aligned}
d\left(H_{\star}([0,1] \otimes T)\right) & =H_{\star}\left(\delta_{0} \otimes T-\delta_{1} \otimes T+[0,1] \otimes d T\right) \\
& =F_{\star} T-G_{\star} T+H_{\star}([0,1] \otimes d T)
\end{aligned}
$$

Therefore we obtain the homotopy formula

$$
\begin{equation*}
F_{\star} T-G_{\star} T=d\left(H_{\star}([0,1] \otimes T)\right)-H_{\star}([0,1] \otimes d T) \tag{2.22}
\end{equation*}
$$

When $T$ is closed, i.e. $d T=0$, we see that $F_{\star} T$ and $G_{\star} T$ are cohomologous on $M_{2}$, i.e. they differ by an exact current $d S$.
$\S$ 2.D.3. Regularization of Currents. Let $\rho \in \mathscr{C}^{\infty}\left(\mathbb{R}^{m}\right)$ be a function with support in $B(0,1)$, such that $\rho(x)$ depends only on $|x|=\left(\sum\left|x_{i}\right|^{2}\right)^{1 / 2}, \rho \geqslant 0$ and $\int_{\mathbb{R}^{m}} \rho(x) d x=1$. We associate to $\rho$ the family of functions $\left(\rho_{\varepsilon}\right)$ such that

$$
\begin{equation*}
\rho_{\varepsilon}(x)=\frac{1}{\varepsilon^{m}} \rho\left(\frac{x}{\varepsilon}\right), \quad \operatorname{Supp} \rho_{\varepsilon} \subset B(0, \varepsilon), \quad \int_{\mathbb{R}^{m}} \rho_{\varepsilon}(x) d x=1 \tag{2.23}
\end{equation*}
$$

We shall refer to this construction by saying that $\left(\rho_{\varepsilon}\right)$ is a family of smoothing kernels. For every current $T=\sum T_{I} d x_{I}$ on an open subset $\Omega \subset \mathbb{R}^{m}$, the family of smooth forms

$$
T \star \rho_{\varepsilon}=\sum_{I}\left(T_{I} \star \rho_{\varepsilon}\right) d x_{I},
$$

defined on $\Omega_{\varepsilon}=\left\{x \in \mathbb{R}^{m} ; d(x, C \Omega)>\varepsilon\right\}$, converges weakly to $T$ as $\varepsilon$ tends to 0 . Indeed, $\left\langle T \star \rho_{\varepsilon}, f\right\rangle=\left\langle T, \rho_{\varepsilon} \star f\right\rangle$ and $\rho_{\varepsilon} \star f$ converges to $f$ in $\mathscr{D}^{p}(\Omega)$ with respect to all seminorms $p_{K}^{s}$.
§ 2.D.4. Poincaré Lemma for Currents. Let $T \in{ }^{s} \mathscr{D}^{\prime}(\Omega)$ be a closed current on an open set $\Omega \subset \mathbb{R}^{m}$. We first show that $T$ is cohomologous to a smooth form. In fact, let $\psi \in \mathscr{C}^{\infty}\left(\mathbb{R}^{m}\right)$ be a cut-off function such that $\operatorname{Supp} \psi \subset \bar{\Omega}, 0<\psi \leqslant 1$ and $|d \psi| \leqslant 1$ on $\Omega$. For any vector $v \in B(0,1)$ we set

$$
F_{v}(x)=x+\psi(x) v .
$$

Since $x \mapsto \psi(x) v$ is a contraction, $F_{v}$ is a diffeomorphism of $\mathbb{R}^{m}$ which leaves $\complement \Omega$ invariant pointwise, so $F_{v}(\Omega)=\Omega$. This diffeomorphism is homotopic to the identity through the homotopy $H_{v}(t, x)=F_{t v}(x):[0,1] \times \Omega \longrightarrow \Omega$ which is proper for every $v$. Formula (2.22) implies

$$
\left(F_{v}\right)_{\star} T-T=d\left(\left(H_{v}\right)_{\star}([0,1] \otimes T)\right) .
$$

After averaging with a smoothing kernel $\rho_{\varepsilon}(v)$ we get $\Theta-T=d S$ where

$$
\Theta=\int_{B(0, \varepsilon)}\left(F_{v}\right)_{\star} T \rho_{\varepsilon}(v) d v, \quad S=\int_{B(0, \varepsilon)}\left(H_{v}\right)_{\star}([0,1] \otimes T) \rho_{\varepsilon}(v) d v
$$

Then $S$ is a current of the same order $s$ as $T$ and $\Theta$ is smooth. Indeed, for $u \in \mathscr{D}^{p}(\Omega)$ we have

$$
\langle\Theta, u\rangle=\left\langle T, u_{\varepsilon}\right\rangle \quad \text { where } \quad u_{\varepsilon}(x)=\int_{B(0, \varepsilon)} F_{v}^{\star} u(x) \rho_{\varepsilon}(v) d v ;
$$

we can make a change of variable $z=F_{v}(x) \Leftrightarrow v=\psi(x)^{-1}(z-x)$ in the last integral and perform derivatives on $\rho_{\varepsilon}$ to see that each seminorm $p_{K}^{t}\left(u_{\varepsilon}\right)$ is controlled by the sup norm of $u$. Thus $\Theta$ and all its derivatives are currents of order 0 , so $\Theta$ is smooth. Now we have $d \Theta=0$ and by the usual Poincaré lemma (1.22) applied to $\Theta$ we obtain
(2.24) Theorem. Let $\Omega \subset \mathbb{R}^{m}$ be a starshaped open subset and $T \in{ }^{s} \mathscr{D}^{\prime q}(\Omega)$ a current of degree $q \geqslant 1$ and order $s$ such that $d T=0$. There exists a current $S \in{ }^{s} \mathscr{D}^{\prime q-1}(\Omega)$ of degree $q-1$ and order $\leqslant s$ such that $d S=T$ on $\Omega$.

## § 3. Holomorphic Functions and Complex Manifolds

## § 3.A. Cauchy Formula in One Variable

We start by recalling a few elementary facts in one complex variable theory. Let $\Omega \subset \mathbb{C}$ be an open set and let $z=x+\mathrm{i} y$ be the complex variable, where $x, y \in \mathbb{R}$. If $f$ is a function of class $C^{1}$ on $\Omega$, we have

$$
d f=\frac{\partial f}{\partial x} d x+\frac{\partial f}{\partial y} d y=\frac{\partial f}{\partial z} d z+\frac{\partial f}{\partial \bar{z}} d \bar{z}
$$

with the usual notations

$$
\begin{equation*}
\frac{\partial}{\partial z}=\frac{1}{2}\left(\frac{\partial}{\partial x}-\mathrm{i} \frac{\partial}{\partial y}\right), \quad \frac{\partial}{\partial \bar{z}}=\frac{1}{2}\left(\frac{\partial}{\partial x}+\mathrm{i} \frac{\partial}{\partial y}\right) . \tag{3.1}
\end{equation*}
$$

The function $f$ is holomorphic on $\Omega$ if $d f$ is $\mathbb{C}$-linear, that is, $\partial f / \partial \bar{z}=0$.
(3.2) Cauchy formula. Let $K \subset \mathbb{C}$ be a compact set with piecewise $C^{1}$ boundary $\partial K$. Then for every $f \in C^{1}(K, \mathbb{C})$

$$
f(w)=\frac{1}{2 \pi \mathrm{i}} \int_{\partial K} \frac{f(z)}{z-w} d z-\int_{K} \frac{1}{\pi(z-w)} \frac{\partial f}{\partial \bar{z}} d \lambda(z), \quad w \in K^{\circ}
$$

where $d \lambda(z)=\frac{i}{2} d z \wedge d \bar{z}=d x \wedge d y$ is the Lebesgue measure on $\mathbb{C}$.
Proof. Assume for simplicity $w=0$. As the function $z \mapsto 1 / z$ is locally integrable at $z=0$, we get

$$
\begin{aligned}
\int_{K} \frac{1}{\pi z} \frac{\partial f}{\partial \bar{z}} d \lambda(z) & =\lim _{\varepsilon \rightarrow 0} \int_{K \backslash D(0, \varepsilon)} \frac{1}{\pi z} \frac{\partial f}{\partial \bar{z}} \frac{\mathrm{i}}{2} d z \wedge d \bar{z} \\
& =\lim _{\varepsilon \rightarrow 0} \int_{K \backslash D(0, \varepsilon)} d\left[\frac{1}{2 \pi \mathrm{i}} f(z) \frac{d z}{z}\right] \\
& =\frac{1}{2 \pi \mathrm{i}} \int_{\partial K} f(z) \frac{d z}{z}-\lim _{\varepsilon \rightarrow 0} \frac{1}{2 \pi \mathrm{i}} \int_{\partial D(0, \varepsilon)} f(z) \frac{d z}{z}
\end{aligned}
$$

by Stokes' formula. The last integral is equal to $\frac{1}{2 \pi} \int_{0}^{2 \pi} f\left(\varepsilon e^{\mathrm{i} \theta}\right) d \theta$ and converges to $f(0)$ as $\varepsilon$ tends to 0 .

When $f$ is holomorphic on $\Omega$, we get the usual Cauchy formula

$$
\begin{equation*}
f(w)=\frac{1}{2 \pi \mathrm{i}} \int_{\partial K} \frac{f(z)}{z-w} d z, \quad w \in K^{\circ}, \tag{3.3}
\end{equation*}
$$

from which many basic properties of holomorphic functions can be derived: power and Laurent series expansions, Cauchy residue formula, ... Another interesting consequence is:
(3.4) Corollary. The $L_{\text {loc }}^{1}$ function $E(z)=1 / \pi z$ is a fundamental solution of the operator $\partial / \partial \bar{z}$ on $\mathbb{C}$, i.e. $\partial E / \partial \bar{z}=\delta_{0}$ (Dirac measure at 0 ). As a consequence, if $v$ is a distribution with compact support in $\mathbb{C}$, then the convolution $u=(1 / \pi z) \star v$ is a solution of the equation $\partial u / \partial \bar{z}=v$.

Proof. Apply (3.2) with $w=0, f \in \mathscr{D}(\mathbb{C})$ and $K \supset \operatorname{Supp} f$, so that $f=0$ on the boundary $\partial K$ and $f(0)=\langle 1 / \pi z,-\partial f / \partial \bar{z}\rangle$.
(3.5) Remark. It should be observed that this formula cannot be used to solve the equation $\partial u / \partial \bar{z}=v$ when $\operatorname{Supp} v$ is not compact; moreover, if Supp $v$ is compact, a solution $u$ with compact support need not always exist. Indeed, we have a necessary condition

$$
\left\langle v, z^{n}\right\rangle=-\left\langle u, \partial z^{n} / \partial \bar{z}\right\rangle=0
$$

for all integers $n \geqslant 0$. Conversely, when the necessary condition $\left\langle v, z^{n}\right\rangle=0$ is satisfied, the canonical solution $u=(1 / \pi z) \star v$ has compact support: this is easily seen by means of the power series expansion $(w-z)^{-1}=\sum z^{n} w^{-n-1}$, if we suppose that $\operatorname{Supp} v$ is contained in the disk $|z|<R$ and that $|w|>R$.

## § 3.B. Holomorphic Functions of Several Variables

Let $\Omega \subset \mathbb{C}^{n}$ be an open set. A function $f: \Omega \rightarrow \mathbb{C}$ is said to be holomorphic if $f$ is continuous and separately holomorphic with respect to each variable, i.e. $z_{j} \mapsto f\left(\ldots, z_{j}, \ldots\right)$ is holomorphic when $z_{1}, \ldots, z_{j-1}, z_{j+1}, \ldots, z_{n}$ are fixed. The set of holomorphic functions on $\Omega$ is a ring and will be denoted $\mathbb{O}(\Omega)$. We first extend the Cauchy formula to the case of polydisks. The open polydisk $D\left(z_{0}, R\right)$ of center $\left(z_{0,1}, \ldots, z_{0, n}\right)$ and (multi)radius $R=\left(R_{1}, \ldots, R_{n}\right)$ is defined as the product of the disks of center $z_{0, j}$ and radius $R_{j}>0$ in each factor $\mathbb{C}$ :

$$
\begin{equation*}
D\left(z_{0}, R\right)=D\left(z_{0,1}, R_{1}\right) \times \ldots \times D\left(z_{0, n}, R_{n}\right) \subset \mathbb{C}^{n} \tag{3.6}
\end{equation*}
$$

The distinguished boundary of $D\left(z_{0}, R\right)$ is by definition the product of the boundary circles

$$
\begin{equation*}
\Gamma\left(z_{0}, R\right)=\Gamma\left(z_{0,1}, R_{1}\right) \times \ldots \times \Gamma\left(z_{0, n}, R_{n}\right) \tag{3.7}
\end{equation*}
$$

It is important to observe that the distinguished boundary is smaller than the topological boundary $\partial D\left(z_{0}, R\right)=\bigcup_{j}\left\{z \in \bar{D}\left(z_{0}, R\right) ;\left|z_{j}-z_{0, j}\right|=R_{j}\right\}$ when $n \geqslant 2$. By induction on $n$, we easily get the
(3.8) Cauchy formula on polydisks. If $\bar{D}\left(z_{0}, R\right)$ is a closed polydisk contained in $\Omega$ and $f \in \mathscr{O}(\Omega)$, then for all $w \in D\left(z_{0}, R\right)$ we have

$$
f(w)=\frac{1}{(2 \pi \mathrm{i})^{n}} \int_{\Gamma\left(z_{0}, R\right)} \frac{f\left(z_{1}, \ldots, z_{n}\right)}{\left(z_{1}-w_{1}\right) \ldots\left(z_{n}-w_{n}\right)} d z_{1} \ldots d z_{n}
$$

The expansion $\left(z_{j}-w_{j}\right)^{-1}=\sum\left(w_{j}-z_{0, j}\right)^{\alpha_{j}}\left(z_{j}-z_{0, j}\right)^{-\alpha_{j}-1}, \alpha_{j} \in \mathbb{N}, 1 \leqslant j \leqslant n$, shows that $f$ can be expanded as a convergent power series $f(w)=\sum_{\alpha \in \mathbb{N}^{n}} a_{\alpha}\left(w-z_{0}\right)^{\alpha}$ over the polydisk $D\left(z_{0}, R\right)$, with the standard notations $z^{\alpha}=z_{1}^{\alpha_{1}} \ldots z_{n}^{\alpha_{n}}, \alpha!=\alpha_{1}!\ldots \alpha_{n}$ ! and with

$$
\begin{equation*}
a_{\alpha}=\frac{1}{(2 \pi \mathrm{i})^{n}} \int_{\Gamma\left(z_{0}, R\right)} \frac{f\left(z_{1}, \ldots, z_{n}\right) d z_{1} \ldots d z_{n}}{\left(z_{1}-z_{0,1}\right)^{\alpha_{1}+1} \ldots\left(z_{n}-z_{0, n}\right)^{\alpha_{n}+1}}=\frac{f^{(\alpha)}\left(z_{0}\right)}{\alpha!} . \tag{3.9}
\end{equation*}
$$

As a consequence, $f$ is holomorphic over $\Omega$ if and only if $f$ is $\mathbb{C}$-analytic. Arguments similar to the one variable case easily yield the
(3.10) Analytic continuation theorem. If $\Omega$ is connected and if there exists a point $z_{0} \in \Omega$ such that $f^{(\alpha)}\left(z_{0}\right)=0$ for all $\alpha \in \mathbb{N}^{n}$, then $f=0$ on $\Omega$.

Another consequence of (3.9) is the Cauchy inequality

$$
\begin{equation*}
\left|f^{(\alpha)}\left(z_{0}\right)\right| \leqslant \frac{\alpha!}{R^{\alpha}} \sup _{\Gamma\left(z_{0}, R\right)}|f|, \quad \bar{D}\left(z_{0}, R\right) \subset \Omega, \tag{3.11}
\end{equation*}
$$

From this, it follows that every bounded holomorphic function on $\mathbb{C}^{n}$ is constant (Liouville's theorem), and more generally, every holomorphic function $F$ on $\mathbb{C}^{n}$ such that
$|F(z)| \leqslant A(1+|z|)^{B}$ with suitable constants $A, B \geqslant 0$ is in fact a polynomial of total degree $\leqslant B$.

We endow $\mathscr{O}(\Omega)$ with the topology of uniform convergence on compact sets $K \subset \subset \Omega$, that is, the topology induced by $C^{0}(\Omega, \mathbb{C})$. Then $\mathscr{O}(\Omega)$ is closed in $C^{0}(\Omega, \mathbb{C})$. The Cauchy inequalities (3.11) show that all derivations $D^{\alpha}$ are continuous operators on $\mathscr{O}(\Omega)$ and that any sequence $f_{j} \in \mathscr{O}(\Omega)$ that is uniformly bounded on all compact sets $K \subset \subset \Omega$ is locally equicontinuous. By Ascoli's theorem, we obtain
(3.12) Montel's theorem. Every locally uniformly bounded sequence $\left(f_{j}\right)$ in $\mathfrak{G}(\Omega)$ has a convergent subsequence $\left(f_{j(\nu)}\right)$.

In other words, bounded subsets of the Fréchet space $\mathfrak{O}(\Omega)$ are relatively compact (a Fréchet space possessing this property is called a Montel space).

## § 3.C. Differential Calculus on Complex Analytic Manifolds

A complex analytic manifold $X$ of dimension $\operatorname{dim}_{\mathbb{C}} X=n$ is a differentiable manifold equipped with a holomorphic atlas $\left(\tau_{\alpha}\right)$ with values in $\mathbb{C}^{n}$; this means by definition that the transition maps $\tau_{\alpha \beta}$ are holomorphic. The tangent spaces $T_{X, x}$ then have a natural complex vector space structure, given by the coordinate isomorphisms

$$
d \tau_{\alpha}(x): T_{X, x} \longrightarrow \mathbb{C}^{n}, \quad U_{\alpha} \ni x
$$

the induced complex structure on $T_{X, x}$ is indeed independent of $\alpha$ since the differentials $d \tau_{\alpha \beta}$ are $\mathbb{C}$-linear isomorphisms. We denote by $T_{X}^{\mathbb{R}}$ the underlying real tangent space and by $J \in \operatorname{End}\left(T_{X}^{\mathbb{R}}\right)$ the almost complex structure, i.e. the operator of multiplication by $\mathrm{i}=\sqrt{-1}$. If $\left(z_{1}, \ldots, z_{n}\right)$ are complex analytic coordinates on an open subset $\Omega \subset X$ and $z_{k}=x_{k}+\mathrm{i} y_{k}$, then $\left(x_{1}, y_{1}, \ldots, x_{n}, y_{n}\right)$ define real coordinates on $\Omega$, and $T_{X \mid \Omega}^{\mathbb{R}}$ admits $\left(\partial / \partial x_{1}, \partial / \partial y_{1}, \ldots, \partial / \partial x_{n}, \partial / \partial y_{n}\right)$ as a basis; the almost complex structure is given by $J\left(\partial / \partial x_{k}\right)=\partial / \partial y_{k}, J\left(\partial / \partial y_{k}\right)=-\partial / \partial x_{k}$. The complexified tangent space $\mathbb{C} \otimes T_{X}=\mathbb{C} \otimes_{\mathbb{R}} T_{X}^{\mathbb{R}}=T_{X}^{\mathbb{R}} \oplus \mathrm{i} T_{X}^{\mathbb{R}}$ splits into conjugate complex subspaces which are the eigenspaces of the complexified endomorphism $\operatorname{Id} \otimes J$ associated to the eigenvalues i and -i . These subspaces have respective bases

$$
\begin{equation*}
\frac{\partial}{\partial z_{k}}=\frac{1}{2}\left(\frac{\partial}{\partial x_{k}}-\mathrm{i} \frac{\partial}{\partial y_{k}}\right), \quad \frac{\partial}{\partial \bar{z}_{k}}=\frac{1}{2}\left(\frac{\partial}{\partial x_{k}}+\mathrm{i} \frac{\partial}{\partial y_{k}}\right), \quad 1 \leqslant k \leqslant n \tag{3.13}
\end{equation*}
$$

and are denoted $T^{1,0} X$ (holomorphic vectors or vectors of type $(1,0)$ ) and $T^{0,1} X$ (antiholomorphic vectors or vectors of type $(0,1))$. The subspaces $T^{1,0} X$ and $T^{0,1} X$ are canonically isomorphic to the complex tangent space $T_{X}$ (with complex structure $J$ ) and its conjugate $\overline{T_{X}}$ (with conjugate complex structure $-J$ ), via the $\mathbb{C}$-linear embeddings

$$
\begin{aligned}
T_{X} & \longrightarrow T_{X}^{1,0} \subset \mathbb{C} \otimes T_{X}, \overline{T_{X}} \longrightarrow T_{X}^{0,1} \subset \mathbb{C} \otimes T_{X} \\
\xi \longmapsto \frac{1}{2}(\xi-\mathrm{i} J \xi), \quad \xi & \longmapsto \frac{1}{2}(\xi+\mathrm{i} J \xi) .
\end{aligned}
$$

We thus have a canonical decomposition $\mathbb{C} \otimes T_{X}=T_{X}^{1,0} \oplus T_{X}^{0,1} \simeq T_{X} \oplus \overline{T_{X}}$, and by duality a decomposition

$$
\operatorname{Hom}_{\mathbb{R}}\left(T_{X}^{\mathbb{R}} ; \mathbb{C}\right) \simeq \operatorname{Hom}_{\mathbb{C}}\left(\mathbb{C} \otimes T_{X} ; \mathbb{C}\right) \simeq T_{X}^{\star} \oplus \overline{T_{X}^{\star}}
$$

where $T_{X}^{\star}$ is the space of $\mathbb{C}$-linear forms and $\overline{T_{X}^{\star}}$ the space of conjugate $\mathbb{C}$-linear forms. With these notations, $\left(d x_{k}, d y_{k}\right)$ is a basis of $\operatorname{Hom}_{\mathbb{R}}\left(T_{\mathbb{R}} X, \mathbb{C}\right),\left(d z_{j}\right)$ a basis of $T_{X}^{\star},\left(d \bar{z}_{j}\right)$ a basis of $\overline{T_{X}^{\star}}$, and the differential of a function $f \in C^{1}(\Omega, \mathbb{C})$ can be written

$$
\begin{equation*}
d f=\sum_{k=1}^{n} \frac{\partial f}{\partial x_{k}} d x_{k}+\frac{\partial f}{\partial y_{k}} d y_{k}=\sum_{k=1}^{n} \frac{\partial f}{\partial z_{k}} d z_{k}+\frac{\partial f}{\partial \bar{z}_{k}} d \bar{z}_{k} . \tag{3.14}
\end{equation*}
$$

The function $f$ is holomorphic on $\Omega$ if and only if $d f$ is $\mathbb{C}$-linear, i.e. if and only if $f$ satisfies the Cauchy-Riemann equations $\partial f / \partial \bar{z}_{k}=0$ on $\Omega, 1 \leqslant k \leqslant n$. We still denote here by $\mathfrak{O}(X)$ the algebra of holomorphic functions on $X$.

Now, we study the basic rules of complex differential calculus. The complexified exterior algebra $\mathbb{C} \otimes_{\mathbb{R}} \Lambda_{\mathbb{R}}^{\bullet}\left(T_{X}^{\mathbb{R}}\right)^{\star}=\Lambda_{\mathbb{C}}^{\bullet}\left(\mathbb{C} \otimes T_{X}\right)^{\star}$ is given by

$$
\Lambda^{k}\left(\mathbb{C} \otimes T_{X}\right)^{\star}=\Lambda^{k}\left(T_{X} \oplus \overline{T_{X}}\right)^{\star}=\bigoplus_{p+q=k} \Lambda^{p, q} T_{X}^{\star}, \quad 0 \leqslant k \leqslant 2 n
$$

where the exterior products are taken over $\mathbb{C}$, and where the components $\Lambda^{p, q} T_{X}^{\star}$ are defined by

$$
\begin{equation*}
\Lambda^{p, q} T_{X}^{\star}=\Lambda^{p} T_{X}^{\star} \otimes \Lambda^{q} \overline{T_{X}^{\star}} . \tag{3.15}
\end{equation*}
$$

A complex differential form $u$ on $X$ is said to be of bidegree or type $(p, q)$ if its value at every point lies in the component $\Lambda^{p, q} T_{X}^{\star}$; we shall denote by $C^{s}\left(\Omega, \Lambda^{p, q} T_{X}^{\star}\right)$ the space of differential forms of bidegree $(p, q)$ and class $C^{s}$ on any open subset $\Omega$ of $X$. If $\Omega$ is a coordinate open set, such a form can be written

$$
u(z)=\sum_{|I|=p,|J|=q} u_{I, J}(z) d z_{I} \wedge d \bar{z}_{J}, \quad u_{I, J} \in C^{s}(\Omega, \mathbb{C}) .
$$

This writing is usually much more convenient than the expression in terms of the real basis $\left(d x_{I} \wedge d y_{J}\right)_{|I|+|J|=k}$ which is not compatible with the splitting of $\Lambda^{k} T_{\mathbb{C}}^{\star} X$ in its $(p, q)$ components. Formula (3.14) shows that the exterior derivative $d$ splits into $d=d^{\prime}+d^{\prime \prime}$, where

$$
\begin{align*}
d^{\prime} & : \mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star}\right) \longrightarrow \mathscr{C}^{\infty}\left(X, \Lambda^{p+1, q} T_{X}^{\star}\right), \\
d^{\prime \prime} & : \mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star}\right) \longrightarrow \mathscr{C}^{\infty}\left(X, \Lambda^{p, q+1} T_{X}^{\star}\right), \\
d^{\prime} u & =\sum_{I, J} \sum_{1 \leqslant k \leqslant n} \frac{\partial u_{I, J}}{\partial z_{k}} d z_{k} \wedge d z_{I} \wedge d \bar{z}_{J}, \\
d^{\prime \prime} u & =\sum_{I, J} \sum_{1 \leqslant k \leqslant n .} \frac{\partial u_{I, J}}{\partial \bar{z}_{k}} d \bar{z}_{k} \wedge d z_{I} \wedge d \bar{z}_{J} .
\end{align*}
$$

The identity $d^{2}=\left(d^{\prime}+d^{\prime \prime}\right)^{2}=0$ is equivalent to

$$
\begin{equation*}
d^{\prime 2}=0, \quad d^{\prime} d^{\prime \prime}+d^{\prime \prime} d^{\prime}=0, \quad d^{\prime 2}=0 \tag{3.17}
\end{equation*}
$$

since these three operators send $(p, q)$-forms in $(p+2, q),(p+1, q+1)$ and $(p, q+2)$-forms, respectively. In particular, the operator $d^{\prime \prime}$ defines for each $p=0,1, \ldots, n$ a complex, called the Dolbeault complex

$$
\mathscr{C}^{\infty}\left(X, \Lambda^{p, 0} T_{X}^{\star}\right) \xrightarrow{d^{\prime \prime}} \cdots \longrightarrow \mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star}\right) \xrightarrow{d^{\prime \prime}} \mathscr{C}^{\infty}\left(X, \Lambda^{p, q+1} T_{X}^{\star}\right)
$$

and corresponding Dolbeault cohomology groups

$$
\begin{equation*}
H^{p, q}(X, \mathbb{C})=\frac{\operatorname{Ker} d^{\prime \prime p, q}}{\operatorname{Im} d^{\prime \prime p, q-1}} \tag{3.18}
\end{equation*}
$$

with the convention that the image of $d^{\prime \prime}$ is zero for $q=0$. The cohomology group $H^{p, 0}(X, \mathbb{C})$ consists of $(p, 0)$-forms $u=\sum_{|I|=p} u_{I}(z) d z_{I}$ such that $\partial u_{I} / \partial \bar{z}_{k}=0$ for all $I, k$, i.e. such that all coefficients $u_{I}$ are holomorphic. Such a form is called a holomorphic $p$-form on $X$.

Let $F: X_{1} \longrightarrow X_{2}$ be a holomorphic map between complex manifolds. The pullback $F^{\star} u$ of a $(p, q)$-form $u$ of bidegree $(p, q)$ on $X_{2}$ is again homogeneous of bidegree $(p, q)$, because the components $F_{k}$ of $F$ in any coordinate chart are holomorphic, hence $F^{\star} d z_{k}=d F_{k}$ is $\mathbb{C}$-linear. In particular, the equality $d F^{\star} u=F^{\star} d u$ implies

$$
\begin{equation*}
d^{\prime} F^{\star} u=F^{\star} d^{\prime} u, \quad d^{\prime \prime} F^{\star} u=F^{\star} d^{\prime \prime} u \tag{3.19}
\end{equation*}
$$

Note that these commutation relations are no longer true for a non holomorphic change of variable. As in the case of the De Rham cohomology groups, we get a pull-back morphism

$$
F^{\star}: H^{p, q}\left(X_{2}, \mathbb{C}\right) \longrightarrow H^{p, q}\left(X_{1}, \mathbb{C}\right) .
$$

The rules of complex differential calculus can be easily extended to currents. We use the following notation.
(3.20) Definition. There are decompositions

$$
\mathscr{D}^{k}(X, \mathbb{C})=\bigoplus_{p+q=k} \mathscr{D}^{p, q}(X, \mathbb{C}), \quad \mathscr{D}_{k}^{\prime}(X, \mathbb{C})=\bigoplus_{p+q=k} \mathscr{D}_{p, q}^{\prime}(X, \mathbb{C}) .
$$

The space $\mathscr{D}_{p, q}^{\prime}(X, \mathbb{C})$ is called the space of currents of bidimension $(p, q)$ and bidegree $(n-p, n-q)$ on $X$, and is also denoted $\mathscr{D}^{\prime n-p, n-q}(X, \mathbb{C})$.

## § 3.D. Newton and Bochner-Martinelli Kernels

The Newton kernel is the elementary solution of the usual Laplace operator $\Delta=$ $\sum \partial^{2} / \partial x_{j}^{2}$ in $\mathbb{R}^{m}$. We first recall a construction of the Newton kernel.

Let $d \lambda=d x_{1} \ldots d x_{m}$ be the Lebesgue measure on $\mathbb{R}^{m}$. We denote by $B(a, r)$ the euclidean open ball of center $a$ and radius $r$ in $\mathbb{R}^{m}$ and by $S(a, r)=\partial B(a, r)$ the corresponding sphere. Finally, we set $\alpha_{m}=\operatorname{Vol}(B(0,1))$ and $\sigma_{m-1}=m \alpha_{m}$ so that

$$
\begin{equation*}
\operatorname{Vol}(B(a, r))=\alpha_{m} r^{m}, \quad \operatorname{Area}(S(a, r))=\sigma_{m-1} r^{m-1} \tag{3.21}
\end{equation*}
$$

The second equality follows from the first by derivation. An explicit computation of the integral $\int_{\mathbb{R}^{m}} e^{-|x|^{2}} d \lambda(x)$ in polar coordinates shows that $\alpha_{m}=\pi^{m / 2} /(m / 2)$ ! where $x!=\Gamma(x+1)$ is the Euler Gamma function. The Newton kernel is then given by:

$$
\begin{cases}N(x)=\frac{1}{2 \pi} \log |x| & \text { if } m=2  \tag{3.22}\\ N(x)=-\frac{1}{(m-2) \sigma_{m-1}}|x|^{2-m} & \text { if } m \neq 2\end{cases}
$$

The function $N(x)$ is locally integrable on $\mathbb{R}^{m}$ and satisfies $\Delta N=\delta_{0}$. When $m=2$, this follows from Cor. 3.4 and the fact that $\Delta=4 \partial^{2} / \partial z \partial \bar{z}$. When $m \neq 2$, this can be checked by computing the weak limit

$$
\begin{aligned}
\lim _{\varepsilon \rightarrow 0} \Delta\left(|x|^{2}+\varepsilon^{2}\right)^{1-m / 2} & =\lim _{\varepsilon \rightarrow 0} m(2-m) \varepsilon^{2}\left(|x|^{2}+\varepsilon^{2}\right)^{-1-m / 2} \\
& =m(2-m) I_{m} \delta_{0}
\end{aligned}
$$

with $I_{m}=\int_{\mathbb{R}^{m}}\left(|x|^{2}+1\right)^{-1-m / 2} d \lambda(x)$. The last equality is easily seen by performing the change of variable $y=\varepsilon x$ in the integral

$$
\int_{\mathbb{R}^{m}} \varepsilon^{2}\left(|x|^{2}+\varepsilon^{2}\right)^{-1-m / 2} f(x) d \lambda(x)=\int_{\mathbb{R}^{m}}\left(|y|^{2}+1\right)^{-1-m / 2} f(\varepsilon y) d \lambda(y)
$$

where $f$ is an arbitrary test function. Using polar coordinates, we find that $I_{m}=\sigma_{m-1} / m$ and our formula follows.

The Bochner-Martinelli kernel is the $(n, n-1)$-differential form on $\mathbb{C}^{n}$ with $L_{\text {loc }}^{1}$ coefficients defined by

$$
\begin{align*}
k_{\mathrm{BM}}(z) & =c_{n} \sum_{1 \leqslant j \leqslant n}(-1)^{j} \frac{\bar{z}_{j} d z_{1} \wedge \ldots d z_{n} \wedge d \bar{z}_{1} \wedge \ldots \widehat{d \bar{z}_{j}} \ldots \wedge d \bar{z}_{n}}{|z|^{2 n}},  \tag{3.23}\\
c_{n} & =(-1)^{n(n-1) / 2} \frac{(n-1)!}{(2 \pi \mathrm{i})^{n}} .
\end{align*}
$$

(3.24) Lemma. $d^{\prime \prime} k_{\mathrm{BM}}=\delta_{0}$ on $\mathbb{C}^{n}$.

Proof. Since the Lebesgue measure on $\mathbb{C}^{n}$ is

$$
d \lambda(z)=\bigwedge_{1 \leqslant j \leqslant n} \frac{\mathrm{i}}{2} d z_{j} \wedge d \bar{z}_{j}=\left(\frac{\mathrm{i}}{2}\right)^{n}(-1)^{\frac{n(n-1)}{2}} d z_{1} \wedge \ldots d z_{n} \wedge d \bar{z}_{1} \wedge \ldots d \bar{z}_{n}
$$

we find

$$
\begin{aligned}
d^{\prime \prime} k_{\mathrm{BM}} & =-\frac{(n-1)!}{\pi^{n}} \sum_{1 \leqslant j \leqslant n} \frac{\partial}{\partial \bar{z}_{j}}\left(\frac{\bar{z}_{j}}{|z|^{2 n}}\right) d \lambda(z) \\
& =-\frac{1}{n(n-1) \alpha_{2 n}} \sum_{1 \leqslant j \leqslant n} \frac{\partial^{2}}{\partial z_{j} \partial \bar{z}_{j}}\left(\frac{1}{|z|^{2 n-2}}\right) d \lambda(z) \\
& =\Delta N(z) d \lambda(z)=\delta_{0} .
\end{aligned}
$$

We let $K_{\mathrm{BM}}(z, \zeta)$ be the pull-back of $k_{\mathrm{BM}}$ by the map $\pi: \mathbb{C}^{n} \times \mathbb{C}^{n} \rightarrow \mathbb{C}^{n},(z, \zeta) \longmapsto$ $z-\zeta$. Then Formula (2.19) implies

$$
\begin{equation*}
d^{\prime \prime} K_{\mathrm{BM}}=\pi^{\star} \delta_{0}=[\Delta], \tag{3.25}
\end{equation*}
$$

where [ $\Delta$ ] denotes the current of integration on the diagonal $\Delta \subset \mathbb{C}^{n} \times \mathbb{C}^{n}$.
(3.26) Koppelman formula. Let $\Omega \subset \mathbb{C}^{n}$ be a bounded open set with piecewise $C^{1}$ boundary. Then for every $(p, q)$-form $v$ of class $C^{1}$ on $\bar{\Omega}$ we have

$$
\begin{aligned}
& v(z)=\int_{\partial \Omega} K_{\mathrm{BM}}^{p, q}(z, \zeta) \wedge v(\zeta) \\
&+d_{z}^{\prime \prime} \int_{\Omega} K_{\mathrm{BM}}^{p, q-1}(z, \zeta) \wedge v(\zeta)+\int_{\Omega} K_{\mathrm{BM}}^{p, q}(z, \zeta) \wedge d^{\prime \prime} v(\zeta)
\end{aligned}
$$

on $\Omega$, where $K_{\mathrm{BM}}^{p, q}(z, \zeta)$ denotes the component of $K_{\mathrm{BM}}(z, \zeta)$ of type $(p, q)$ in $z$ and ( $n-$ $p, n-q-1)$ in $\zeta$.
Proof. Given $w \in \mathscr{D}^{n-p, n-q}(\Omega)$, we consider the integral

$$
\int_{\partial \Omega \times \Omega} K_{\mathrm{BM}}(z, \zeta) \wedge v(\zeta) \wedge w(z)
$$

It is well defined since $K_{\mathrm{BM}}$ has no singularities on $\partial \Omega \times \operatorname{Supp} v \subset \subset \partial \Omega \times \Omega$. Since $w(z)$ vanishes on $\partial \Omega$ the integral can be extended as well to $\partial(\Omega \times \Omega)$. As $K_{\mathrm{BM}}(z, \zeta) \wedge v(\zeta) \wedge w(z)$ is of total bidegree $(2 n, 2 n-1)$, its differential $d^{\prime}$ vanishes. Hence Stokes' formula yields

$$
\begin{aligned}
& \int_{\partial \Omega \times \Omega} K_{\mathrm{BM}}(z, \zeta) \wedge v(\zeta) \wedge w(z)=\int_{\Omega \times \Omega} d^{\prime \prime}\left(K_{\mathrm{BM}}(z, \zeta) \wedge v(\zeta) \wedge w(z)\right) \\
& =\int_{\Omega \times \Omega} d^{\prime \prime} K_{\mathrm{BM}}(z, \zeta) \wedge v(\zeta) \wedge w(z)-K_{\mathrm{BM}}^{p, q}(z, \zeta) \wedge d^{\prime \prime} v(\zeta) \wedge w(z) \\
& \quad-(-1)^{p+q} \int_{\Omega \times \Omega} K_{\mathrm{BM}}^{p, q-1}(z, \zeta) \wedge v(\zeta) \wedge d^{\prime \prime} w(z)
\end{aligned}
$$

By (3.25) we have

$$
\int_{\Omega \times \Omega} d^{\prime \prime} K_{\mathrm{BM}}(z, \zeta) \wedge v(\zeta) \wedge w(z)=\int_{\Omega \times \Omega}[\Delta] \wedge v(\zeta) \wedge w(z)=\int_{\Omega} v(z) \wedge w(z)
$$

Denoting $\langle$,$\rangle the pairing between currents and test forms on \Omega$, the above equality is thus equivalent to

$$
\begin{aligned}
\left\langle\int_{\partial \Omega} K_{\mathrm{BM}}(z, \zeta) \wedge v(\zeta), w(z)\right\rangle & =\left\langle v(z)-\int_{\Omega} K_{\mathrm{BM}}^{p, q}(z, \zeta) \wedge d^{\prime \prime} v(\zeta), w(z)\right\rangle \\
& -(-1)^{p+q}\left\langle\int_{\Omega} K_{\mathrm{BM}}^{p, q-1}(z, \zeta) \wedge v(\zeta), d^{\prime \prime} w(z)\right\rangle
\end{aligned}
$$

which is itself equivalent to the Koppelman formula by integrating $d^{\prime \prime} v$ by parts.
(3.27) Corollary. Let $v \in{ }^{s} \mathscr{D}^{p, q}\left(\mathbb{C}^{n}\right)$ be a form of class $C^{s}$ with compact support such that $d^{\prime \prime} v=0, q \geqslant 1$. Then the $(p, q-1)$-form

$$
u(z)=\int_{\mathbb{C}^{n}} K_{\mathrm{BM}}^{p, q-1}(z, \zeta) \wedge v(\zeta)
$$

is a $C^{s}$ solution of the equation $d^{\prime \prime} u=v$. Moreover, if $(p, q)=(0,1)$ and $n \geqslant 2$ then $u$ has compact support, thus the Dolbeault cohomology group with compact support $H_{c}^{0,1}\left(\mathbb{C}^{n}, \mathbb{C}\right)$ vanishes for $n \geqslant 2$.

Proof. Apply the Koppelman formula on a sufficiently large ball $\bar{\Omega}=\bar{B}(0, R)$ containing $\operatorname{Supp} v$. Then the formula immediately gives $d^{\prime \prime} u=v$. Observe that the coefficients of $K_{\mathrm{BM}}(z, \zeta)$ are $O\left(|z-\zeta|^{-(2 n-1)}\right)$, hence $|u(z)|=O\left(|z|^{-(2 n-1)}\right)$ at infinity. If $q=1$, then $u$ is holomorphic on $\mathbb{C}^{n} \backslash \bar{B}(0, R)$. Now, this complement is a union of complex lines when $n \geqslant 2$, hence $u=0$ on $\mathbb{C}^{n} \backslash \bar{B}(0, R)$ by Liouville's theorem.
(3.28) Hartogs extension theorem. Let $\Omega$ be an open set in $\mathbb{C}^{n}, n \geqslant 2$, and let $K \subset \Omega$ be a compact subset such that $\Omega \backslash K$ is connected. Then every holomorphic function $f \in \mathscr{O}(\Omega \backslash K)$ extends into a function $\widetilde{f} \in \mathscr{O}(\Omega)$.
Proof. Let $\psi \in \mathscr{D}(\Omega)$ be a cut-off function equal to 1 on a neighborhood of $K$. Set $f_{0}=(1-\psi) f \in \mathscr{C}^{\infty}(\Omega)$, defined as 0 on $K$. Then $v=d^{\prime \prime} f_{0}=-f d^{\prime \prime} \psi$ can be extended by 0 outside $\Omega$, and can thus be seen as a smooth ( 0,1 )-form with compact support in $\mathbb{C}^{n}$, such that $d^{\prime \prime} v=0$. By Cor. 3.27, there is a smooth function $u$ with compact support in $\mathbb{C}^{n}$ such that $d^{\prime \prime} u=v$. Then $\widetilde{f}=f_{0}-u \in \mathscr{O}(\Omega)$. Now $u$ is holomorphic outside $\operatorname{Supp} \psi$, so $u$ vanishes on the unbounded component $G$ of $\mathbb{C}^{n} \backslash \operatorname{Supp} \psi$. The boundary $\partial G$ is contained in $\partial \operatorname{Supp} \psi \subset \Omega \backslash K$, so $\widetilde{f}=(1-\psi) f-u$ coincides with $f$ on the non empty open set $\Omega \cap G \subset \Omega \backslash K$. Therefore $\tilde{f}=f$ on the connected open set $\Omega \backslash K$.

A refined version of the Hartogs extension theorem due to Bochner will be given in Exercise 8.13. It shows that $f$ need only be given as a $C^{1}$ function on $\partial \Omega$, satisfying the tangential Cauchy-Riemann equations (a so-called CR-function). Then $f$ extends as a holomorphic function $\widetilde{f} \in \mathscr{O}(\Omega) \cap C^{0}(\bar{\Omega})$, provided that $\partial \Omega$ is connected.

## § 3.E. The Dolbeault-Grothendieck Lemma

We are now in a position to prove the Dolbeault-Grothendieck lemma [Dolbeault 1953], which is the analogue for $d^{\prime \prime}$ of the Poincaré lemma. The proof given below makes use of the Bochner-Martinelli kernel. Many other proofs can be given, e.g. by using a reduction to the one dimensional case in combination with the Cauchy formula (3.2), see Exercise 8.5 or [Hörmander 1966].
(3.29) Dolbeault-Grothendieck lemma. Let $\Omega$ be a neighborhood of 0 in $\mathbb{C}^{n}$ and $v \in{ }^{s} \mathbb{E}^{p, q}(\Omega, \mathbb{C}),\left[\right.$ resp. $\left.v \in{ }^{s} \mathscr{D}^{\prime p, q}(\Omega, \mathbb{C})\right]$, such that $d^{\prime \prime} v=0$, where $1 \leqslant s \leqslant \infty$.
a) If $q=0$, then $v(z)=\sum_{|I|=p} v_{I}(z) d z_{I}$ is a holomorphic $p$-form, i.e. a form whose coefficients are holomorphic functions.
b) If $q \geqslant 1$, there exists a neighborhood $\omega \subset \Omega$ of 0 and a form $u$ in ${ }^{s} \mathscr{C}^{p, q-1}(\omega, \mathbb{C})[$ resp. a current $\left.u \in{ }^{S} \mathscr{D}^{\prime p, q-1}(\omega, \mathbb{C})\right]$ such that $d^{\prime \prime} u=v$ on $\omega$.

Proof. We assume that $\Omega$ is a ball $B(0, r) \subset \mathbb{C}^{n}$ and take for simplicity $r>1$ (possibly after a dilation of coordinates). We then set $\omega=B(0,1)$. Let $\psi \in \mathscr{D}(\Omega)$ be a cut-off function equal to 1 on $\omega$. The Koppelman formula (3.26) applied to the form $\psi v$ on $\Omega$ gives

$$
\psi(z) v(z)=d_{z}^{\prime \prime} \int_{\Omega} K_{\mathrm{BM}}^{p, q-1}(z, \zeta) \wedge \psi(\zeta) v(\zeta)+\int_{\Omega} K_{\mathrm{BM}}^{p, q}(z, \zeta) \wedge d^{\prime \prime} \psi(\zeta) \wedge v(\zeta)
$$

This formula is valid even when $v$ is a current, because we may regularize $v$ as $v \star \rho_{\varepsilon}$ and take the limit. We introduce on $\mathbb{C}^{n} \times \mathbb{C}^{n} \times \mathbb{C}^{n}$ the kernel

$$
K(z, w, \zeta)=c_{n} \sum_{j=1}^{n} \frac{(-1)^{j}\left(w_{j}-\bar{\zeta}_{j}\right)}{((z-\zeta) \cdot(w-\bar{\zeta}))^{n}} \bigwedge_{k}\left(d z_{k}-d \zeta_{k}\right) \wedge \bigwedge_{k \neq j}\left(d w_{k}-d \bar{\zeta}_{k}\right)
$$

By construction, $K_{\mathrm{Bm}}(z, \zeta)$ is the result of the substitution $w=\bar{z}$ in $K(z, w, \zeta)$, i.e. $K_{\mathrm{BM}}=h^{\star} K$ where $h(z, \zeta)=(z, \bar{z}, \zeta)$. We denote by $K^{p, q}$ the component of $K$ of bidegree $(p, 0)$ in $z,(q, 0)$ in $w$ and $(n-p, n-q-1)$ in $\zeta$. Then $K_{\mathrm{BM}}^{p, q}=h^{\star} K^{p, q}$ and we find

$$
v=d^{\prime \prime} u_{0}+g^{\star} v_{1} \quad \text { on } \omega,
$$

where $g(z)=(z, \bar{z})$ and

$$
\begin{aligned}
u_{0}(z) & =\int_{\Omega} K_{\mathrm{BM}}^{p, q-1}(z, \zeta) \wedge \psi(\zeta) v(\zeta), \\
v_{1}(z, w) & =\int_{\Omega} K^{p, q}(z, w, \zeta) \wedge d^{\prime \prime} \psi(\zeta) \wedge v(\zeta)
\end{aligned}
$$

By definition of $K^{p, q}(z, w, \zeta), v_{1}$ is holomorphic on the open set

$$
U=\{(z, w) \in \omega \times \omega ; \forall \zeta \notin \omega, \operatorname{Re}(z-\zeta) \cdot(w-\bar{\zeta})>0\}
$$

which contains the "conjugate-diagonal" points $(z, \bar{z})$ as well as the points $(z, 0)$ and $(0, w)$ in $\omega \times \omega$. Moreover $U$ clearly has convex slices $\left(\{z\} \times \mathbb{C}^{n}\right) \cap U$ and $\left(\mathbb{C}^{n} \times\{w\}\right) \cap U$. In particular $U$ is starshaped with respect to $w$, i.e.

$$
(z, w) \in U \Longrightarrow(z, t w) \in U, \quad \forall t \in[0,1] .
$$

As $u_{1}$ is of type $(p, 0)$ in $z$ and $(q, 0)$ in $w$, we get $d_{z}^{\prime \prime}\left(g^{\star} v_{1}\right)=g^{\star} d_{w} v_{1}=0$, hence $d_{w} v_{1}=0$. For $q=0$ we have $K_{\mathrm{BM}}^{p, q-1}=0$, thus $u_{0}=0$, and $v_{1}$ does not depend on $w$, thus $v$ is holomorphic on $\omega$. For $q \geqslant 1$, we can use the homotopy formula (1.23) with respect to $w$ (considering $z$ as a parameter) to get a holomorphic form $u_{1}(z, w)$ of type $(p, 0)$ in $z$ and $(q-1,0)$ in $w$, such that $d_{w} u_{1}(z, w)=v_{1}(z, w)$. Then we get $d^{\prime \prime} g^{\star} u_{1}=g^{\star} d_{w} u_{1}=g^{\star} v_{1}$, hence

$$
v=d^{\prime \prime}\left(u_{0}+g^{\star} u_{1}\right) \quad \text { on } \omega .
$$

Finally, the coefficients of $u_{0}$ are obtained as linear combinations of convolutions of the coefficients of $\psi v$ with $L_{\text {loc }}^{1}$ functions of the form $\bar{\zeta}_{j}|\zeta|^{-2 n}$. Hence $u_{0}$ is of class $C^{s}$ (resp. is a current of order $s$ ), if $v$ is.
(3.30) Corollary. The operator $d^{\prime \prime}$ is hypoelliptic in bidegree ( $p, 0$ ), i.e. if a current $f \in \mathscr{D}^{\prime p, 0}(X, \mathbb{C})$ satisfies $d^{\prime \prime} f \in \mathscr{E}^{p, 1}(X, \mathbb{C})$, then $f \in \mathscr{E}^{p, 0}(X, \mathbb{C})$.

Proof. The result is local, so we may assume that $X=\Omega$ is a neighborhood of 0 in $\mathbb{C}^{n}$. The $(p, 1)$-form $v=d^{\prime \prime} f \in \mathscr{E}^{p, 1}(X, \mathbb{C})$ satisfies $d^{\prime \prime} v=0$, hence there exists $u \in \mathscr{C}^{p, 0}(\widetilde{\Omega}, \mathbb{C})$ such that $d^{\prime \prime} u=d^{\prime \prime} f$. Then $f-u$ is holomorphic and $f=(f-u)+u \in \mathscr{E}^{p, 0}(\widetilde{\Omega}, \mathbb{C})$.

## § 4. Subharmonic Functions

A harmonic (resp. subharmonic) function on an open subset of $\mathbb{R}^{m}$ is essentially a function (or distribution) $u$ such that $\Delta u=0$ (resp. $\Delta u \geqslant 0$ ). A fundamental example of subharmonic function is given by the Newton kernel $N$, which is actually harmonic on $\mathbb{R}^{m} \backslash\{0\}$. Subharmonic functions are an essential tool of harmonic analysis and potential theory. Before giving their precise definition and properties, we derive a basic integral formula involving the Green kernel of the Laplace operator on the ball.

## $\S 4 . A$. Construction of the Green Kernel

The Green kernel $G_{\Omega}(x, y)$ of a smoothly bounded domain $\Omega \subset \subset \mathbb{R}^{m}$ is the solution of the following Dirichlet boundary problem for the Laplace operator $\Delta$ on $\Omega$ :
(4.1) Definition. The Green kernel of a smoothly bounded domain $\Omega \subset \subset \mathbb{R}^{m}$ is a function $G_{\Omega}(x, y): \bar{\Omega} \times \bar{\Omega} \rightarrow[-\infty, 0]$ with the following properties:
a) $G_{\Omega}(x, y)$ is $\mathscr{C}^{\infty}$ on $\bar{\Omega} \times \bar{\Omega} \backslash \operatorname{Diag}_{\Omega} \quad\left(\operatorname{Diag}_{\Omega}=\right.$ diagonal $)$;
b) $G_{\Omega}(x, y)=G_{\Omega}(y, x)$;
c) $G_{\Omega}(x, y)<0$ on $\Omega \times \Omega$ and $G_{\Omega}(x, y)=0$ on $\partial \Omega \times \Omega$;
d) $\Delta_{x} G_{\Omega}(x, y)=\delta_{y}$ on $\Omega$ for every fixed $y \in \Omega$.

It can be shown that $G_{\Omega}$ always exists and is unique. The uniqueness is an easy consequence of the maximum principle (see Th. 4.14 below). In the case where $\Omega=$ $B(0, r)$ is a ball (the only case we are going to deal with), the existence can be shown through explicit calculations. In fact the Green kernel $G_{r}(x, y)$ of $B(0, r)$ is

$$
\begin{equation*}
G_{r}(x, y)=N(x-y)-N\left(\frac{|y|}{r}\left(x-\frac{r^{2}}{|y|^{2}} y\right)\right), \quad x, y \in \bar{B}(0, r) \tag{4.2}
\end{equation*}
$$

A substitution of the explicit value of $N(x)$ yields:

$$
\begin{aligned}
& G_{r}(x, y)=\frac{1}{4 \pi} \log \frac{|x-y|^{2}}{r^{2}-2\langle x, y\rangle+\frac{1}{r^{2}}|x|^{2}|y|^{2}} \quad \text { if } m=2, \quad \text { otherwise } \\
& G_{r}(x, y)=\frac{-1}{(m-2) \sigma_{m-1}}\left(|x-y|^{2-m}-\left(r^{2}-2\langle x, y\rangle+\frac{1}{r^{2}}|x|^{2}|y|^{2}\right)^{1-m / 2}\right) .
\end{aligned}
$$

(4.3) Theorem. The above defined function $G_{r}$ satisfies all four properties (4.1 a-d) on $\Omega=B(0, r)$, thus $G_{r}$ is the Green kernel of $B(0, r)$.

Proof. The first three properties are immediately verified on the formulas, because

$$
r^{2}-2\langle x, y\rangle+\frac{1}{r^{2}}|x|^{2}|y|^{2}=|x-y|^{2}+\frac{1}{r^{2}}\left(r^{2}-|x|^{2}\right)\left(r^{2}-|y|^{2}\right) .
$$

For property d), observe that $r^{2} y /|y|^{2} \notin \bar{B}(0, r)$ whenever $y \in B(0, r) \backslash\{0\}$. The second Newton kernel in the right hand side of (4.1) is thus harmonic in $x$ on $B(0, r)$, and

$$
\Delta_{x} G_{r}(x, y)=\Delta_{x} N(x-y)=\delta_{y} \quad \text { on } B(0, r) .
$$

## §4.B. Green-Riesz Representation Formula and Dirichlet Problem

§ 4.B.1. Green-Riesz Formula. For all smooth functions $u, v$ on a smoothly bounded domain $\Omega \subset \subset \mathbb{R}^{m}$, we have

$$
\begin{equation*}
\int_{\Omega}(u \Delta v-v \Delta u) d \lambda=\int_{\partial \Omega}\left(u \frac{\partial v}{\partial \nu}-v \frac{\partial u}{\partial \nu}\right) d \sigma \tag{4.4}
\end{equation*}
$$

where $\partial / \partial \nu$ is the derivative along the outward normal unit vector $\nu$ of $\partial \Omega$ and $d \sigma$ the euclidean area measure. Indeed

$$
(-1)^{j-1} d x_{1} \wedge \ldots \wedge \widehat{d x_{j}} \wedge \ldots \wedge d x_{m\lceil } \supset \Omega=\nu_{j} d \sigma
$$

for the wedge product of $\langle\nu, d x\rangle$ with the left hand side is $\nu_{j} d \lambda$. Therefore

$$
\frac{\partial v}{\partial \nu} d \sigma=\sum_{j=1}^{m} \frac{\partial v}{\partial x_{j}} \nu_{j} d \sigma=\sum_{j=1}^{m}(-1)^{j-1} \frac{\partial v}{\partial x_{j}} d x_{1} \wedge \ldots \wedge \widehat{d x_{j}} \wedge \ldots \wedge d x_{m} .
$$

Formula (4.4) is then an easy consequence of Stokes' theorem. Observe that (4.4) is still valid if $v$ is a distribution with singular support relatively compact in $\Omega$. For $\Omega=B(0, r)$, $u \in C^{2}(\bar{B}(0, r), \mathbb{R})$ and $v(y)=G_{r}(x, y)$, we get the Green-Riesz representation formula:

$$
\begin{equation*}
u(x)=\int_{B(0, r)} \Delta u(y) G_{r}(x, y) d \lambda(y)+\int_{S(0, r)} u(y) P_{r}(x, y) d \sigma(y) \tag{4.5}
\end{equation*}
$$

where $P_{r}(x, y)=\partial G_{r}(x, y) / \partial \nu(y),(x, y) \in B(0, r) \times S(0, r)$. The function $P_{r}(x, y)$ is called the Poisson kernel. It is smooth and satisfies $\Delta_{x} P_{r}(x, y)=0$ on $B(0, r)$ by (4.1 d). A simple computation left to the reader yields:

$$
\begin{equation*}
P_{r}(x, y)=\frac{1}{\sigma_{m-1} r} \frac{r^{2}-|x|^{2}}{|x-y|^{m}} . \tag{4.6}
\end{equation*}
$$

Formula (4.5) for $u \equiv 1$ shows that $\int_{S(0, r)} P_{r}(x, y) d \sigma(y)=1$. When $x$ in $B(0, r)$ tends to $x_{0} \in S(0, r)$, we see that $P_{r}(x, y)$ converges uniformly to 0 on every compact subset of $S(0, r) \backslash\left\{x_{0}\right\}$; it follows that the measure $P_{r}(x, y) d \sigma(y)$ converges weakly to $\delta_{x_{0}}$ on $S(0, r)$.
§ 4.B.2. Solution of the Dirichlet Problem. For any bounded measurable function $v$ on $S(a, r)$ we define

$$
\begin{equation*}
P_{a, r}[v](x)=\int_{S(a, r)} v(y) P_{r}(x-a, y-a) d \sigma(y), \quad x \in B(a, r) \tag{4.7}
\end{equation*}
$$

If $u \in C^{0}(\bar{B}(a, r), \mathbb{R}) \cap C^{2}(B(a, r), \mathbb{R})$ is harmonic, i.e. $\Delta u=0$ on $B(a, r)$, then (4.5) gives $u=P_{a, r}[u]$ on $B(a, r)$, i.e. the Poisson kernel reproduces harmonic functions. Suppose now that $v \in C^{0}(S(a, r), \mathbb{R})$ is given. Then $P_{r}(x-a, y-a) d \sigma(y)$ converges weakly to $\delta_{x_{0}}$ when $x$ tends to $x_{0} \in S(a, r)$, so $P_{a, r}[v](x)$ converges to $v\left(x_{0}\right)$. It follows that the function $u$ defined by

$$
\left\{\begin{array}{lll}
u=P_{a, r}[v] & \text { on } & B(a, r), \\
u=v & \text { on } & S(a, r)
\end{array}\right.
$$

is continuous on $\bar{B}(a, r)$ and harmonic on $B(a, r)$; thus $u$ is the solution of the Dirichlet problem with boundary values $v$.

## § 4.C. Definition and Basic Properties of Subharmonic Functions

$\S$ 4.C.1. Definition. Mean Value Inequalities. If $u$ is a Borel function on $\bar{B}(a, r)$ which is bounded above or below, we consider the mean values of $u$ over the ball or sphere:

$$
\begin{align*}
& \mu_{B}(u ; a, r)=\frac{1}{\alpha_{m} r^{m}} \int_{B(a, r)} u(x) d \lambda(x),  \tag{4.8}\\
& \mu_{S}(u ; a, r)=\frac{1}{\sigma_{m-1} r^{m-1}} \int_{S(a, r)} u(x) d \sigma(x) .
\end{align*}
$$

As $d \lambda=d r d \sigma$ these mean values are related by

$$
\begin{align*}
\mu_{B}(u ; a, r) & =\frac{1}{\alpha_{m} r^{m}} \int_{0}^{r} \sigma_{m-1} t^{m-1} \mu_{S}(u ; a, t) d t  \tag{4.9}\\
& =m \int_{0}^{1} t^{m-1} \mu_{S}(u ; a, r t) d t
\end{align*}
$$

Now, apply formula (4.5) with $x=0$. We get $P_{r}(0, y)=1 / \sigma_{m-1} r^{m-1}$ and $G_{r}(0, y)=$ $\left(|y|^{2-m}-r^{2-m}\right) /(2-m) \sigma_{m-1}=-\left(1 / \sigma_{m-1}\right) \int_{|y|}^{r} t^{1-m} d t$, thus

$$
\begin{aligned}
\int_{B(0, r)} \Delta u(y) G_{r}(0, y) d \lambda(y) & =-\frac{1}{\sigma_{m-1}} \int_{0}^{r} \frac{d t}{t^{m-1}} \int_{|y|<t} \Delta u(y) d \lambda(y) \\
& =-\frac{1}{m} \int_{0}^{r} \mu_{B}(\Delta u ; 0, t) t d t
\end{aligned}
$$

thanks to the Fubini formula. By translating $S(0, r)$ to $S(a, r)$, (4.5) implies the Gauss formula

$$
\begin{equation*}
\mu_{S}(u ; a, r)=u(a)+\frac{1}{m} \int_{0}^{r} \mu_{B}(\Delta u ; a, t) t d t \tag{4.10}
\end{equation*}
$$

Let $\Omega$ be an open subset of $\mathbb{R}^{m}$ and $u \in C^{2}(\Omega, \mathbb{R})$. If $a \in \Omega$ and $\Delta u(a)>0$ (resp. $\Delta u(a)<0$ ), Formula (4.10) shows that $\mu_{S}(u ; a, r)>u(a)$ (resp. $\left.\mu_{S}(u ; a, r)<u(a)\right)$ for $r$ small enough. In particular, $u$ is harmonic (i.e. $\Delta u=0$ ) if and only if $u$ satisfies the mean value equality

$$
\mu_{S}(u ; a, r)=u(a), \quad \forall \bar{B}(a, r) \subset \Omega .
$$

Now, observe that if $\left(\rho_{\varepsilon}\right)$ is a family of radially symmetric smoothing kernels associated with $\rho(x)=\widetilde{\rho}(|x|)$ and if $u$ is a Borel locally bounded function, an easy computation yields

$$
\begin{align*}
u \star \rho_{\varepsilon}(a) & =\int_{B(0,1)} u(a+\varepsilon x) \rho(x) d \lambda \\
& =\sigma_{m-1} \int_{0}^{1} \mu_{S}(u ; a, \varepsilon t) \widetilde{\rho}(t) t^{m-1} d t \tag{4.11}
\end{align*}
$$

Thus, if $u$ is a Borel locally bounded function satisfying the mean value equality on $\Omega$, (4.11) shows that $u \star \rho_{\varepsilon}=u$ on $\Omega_{\varepsilon}$, in particular $u$ must be smooth. Similarly, if we replace the mean value equality by an inequality, the relevant regularity property to be required for $u$ is just semicontinuity.
(4.12) Theorem and definition. Let $u: \Omega \longrightarrow[-\infty,+\infty[$ be an upper semicontinuous function. The following various forms of mean value inequalities are equivalent:
a) $u(x) \leqslant P_{a, r}[u](x), \quad \forall \bar{B}(a, r) \subset \Omega, \quad \forall x \in B(a, r)$;
b) $u(a) \leqslant \mu_{S}(u ; a, r), \quad \forall \bar{B}(a, r) \subset \Omega$;
c) $u(a) \leqslant \mu_{B}(u ; a, r), \quad \forall \bar{B}(a, r) \subset \Omega$;
d) for every $a \in \Omega$, there exists a sequence $\left(r_{\nu}\right)$ decreasing to 0 such that

$$
u(a) \leqslant \mu_{B}\left(u ; a, r_{\nu}\right) \quad \forall \nu ;
$$

e) for every $a \in \Omega$, there exists a sequence $\left(r_{\nu}\right)$ decreasing to 0 such that

$$
u(a) \leqslant \mu_{S}\left(u ; a, r_{\nu}\right) \quad \forall \nu .
$$

A function $u$ satisfying one of the above properties is said to be subharmonic on $\Omega$. The set of subharmonic functions will be denoted by $\operatorname{Sh}(\Omega)$.

By (4.10) we see that a function $u \in C^{2}(\Omega, \mathbb{R})$ is subharmonic if and only if $\Delta u \geqslant 0$ : in fact $\mu_{S}(u ; a, r)<u(a)$ for $r$ small if $\Delta u(a)<0$. It is also clear on the definitions that every (locally) convex function on $\Omega$ is subharmonic.

Proof. We have obvious implications

$$
\text { a) } \Longrightarrow b) \Longrightarrow c) \Longrightarrow d) \Longrightarrow e),
$$

the second and last ones by (4.10) and the fact that $\mu_{B}\left(u ; a, r_{\nu}\right) \leqslant \mu_{S}(u ; a, t)$ for at least one $t \in] 0, r_{\nu}[$. In order to prove e$) \Longrightarrow \mathrm{a}$ ), we first need a suitable version of the maximum principle.
(4.13) Lemma. Let $u: \Omega \longrightarrow[-\infty,+\infty[$ be an upper semicontinuous function satisfying property 4.12 e ). If $u$ attains its supremum at a point $x_{0} \in \Omega$, then $u$ is constant on the connected component of $x_{0}$ in $\Omega$.

Proof. We may assume that $\Omega$ is connected. Let

$$
W=\left\{x \in \Omega ; u(x)<u\left(x_{0}\right)\right\} .
$$

$W$ is open by the upper semicontinuity, and distinct from $\Omega$ since $x_{0} \notin W$. We want to show that $W=\emptyset$. Otherwise $W$ has a non empty connected component $W_{0}$, and $W_{0}$ has a boundary point $a \in \Omega$. We have $a \in \Omega \backslash W$, thus $u(a)=u\left(x_{0}\right)$. By assumption 4.12 e ), we get $u(a) \leqslant \mu_{S}\left(u ; a, r_{\nu}\right)$ for some sequence $r_{\nu} \rightarrow 0$. For $r_{\nu}$ small enough, $W_{0}$ intersects $\Omega \backslash \bar{B}\left(a, r_{\nu}\right)$ and $B\left(a, r_{\nu}\right)$; as $W_{0}$ is connected, we also have $S\left(a, r_{\nu}\right) \cap W_{0} \neq \emptyset$. Since $u \leqslant u\left(x_{0}\right)$ on the sphere $S\left(a, r_{\nu}\right)$ and $u<u\left(x_{0}\right)$ on its open subset $S\left(a, r_{\nu}\right) \cap W_{0}$, we get $u(a) \leqslant \mu_{S}(u ; a, r)<u\left(x_{0}\right)$, a contradiction.
(4.14) Maximum principle. If $u$ is subharmonic in $\Omega$ (in the sense that $u$ satisfies the weakest property 4.12 e$)$ ), then

$$
\sup _{\Omega} u=\limsup _{\Omega \ni z \rightarrow \partial \Omega \cup\{\infty\}} u(z),
$$

and $\sup _{K} u=\sup _{\partial K} u(z)$ for every compact subset $K \subset \Omega$.
Proof. We have of course $\lim \sup _{z \rightarrow \partial \Omega \cup\{\infty\}} u(z) \leqslant \sup _{\Omega} u$. If the inequality is strict, this means that the supremum is achieved on some compact subset $L \subset \Omega$. Thus, by the upper semicontinuity, there is $x_{0} \in L$ such that $\sup _{\Omega} u=\sup _{L} u=u\left(x_{0}\right)$. Lemma 4.13 shows that $u$ is constant on the connected component $\Omega_{0}$ of $x_{0}$ in $\Omega$, hence

$$
\sup _{\Omega} u=u\left(x_{0}\right)=\limsup _{\Omega_{0} \ni z \rightarrow \partial \Omega_{0} \cup\{\infty\}} u(z) \leqslant \limsup _{\Omega \ni z \rightarrow \partial \Omega \cup\{\infty\}} u(z),
$$

contradiction. The statement involving a compact subset $K$ is obtained by applying the first statement to $\Omega^{\prime}=K^{\circ}$.

Proof of (4.12) e) $\Longrightarrow$ a). Let $u$ be an upper semicontinuous function satisfying 4.12 e) and $\bar{B}(a, r) \subset \Omega$ an arbitrary closed ball. One can find a decreasing sequence of continuous functions $v_{k} \in C^{0}(S(a, r), \mathbb{R})$ such that $\lim v_{k}=u$. Set $h_{k}=P_{a, r}\left[v_{k}\right] \in$ $C^{0}(\bar{B}(a, r), \mathbb{R})$. As $h_{k}$ is harmonic on $B(a, r)$, the function $u-h_{k}$ satisfies 4.12 e) on $B(a, r)$. Furthermore $\limsup _{x \rightarrow \xi \in S(a, r)} u(x)-h_{k}(x) \leqslant u(\xi)-v_{k}(\xi) \leqslant 0$, so $u-h_{k} \leqslant 0$ on $B(a, r)$ by Th. 4.14. By monotone convergence, we find $u \leqslant P_{a, r}[u]$ on $B(a, r)$ when $k$ tends to $+\infty$.
§ 4.C.2. Basic Properties. Here is a short list of the most basic properties.
(4.15) Theorem. For any decreasing sequence $\left(u_{k}\right)$ of subharmonic functions, the limit $u=\lim u_{k}$ is subharmonic.

Proof. A decreasing limit of upper semicontinuous functions is again upper semicontinuous, and the mean value inequalities 4.12 remain valid for $u$ by Lebesgue's monotone convergence theorem.
(4.16) Theorem. Let $u_{1}, \ldots, u_{p} \in \operatorname{Sh}(\Omega)$ and $\chi: \mathbb{R}^{p} \longrightarrow \mathbb{R}$ be a convex function such that $\chi\left(t_{1}, \ldots, t_{p}\right)$ is non decreasing in each $t_{j}$. If $\chi$ is extended by continuity into a function $\left[-\infty,+\infty\left[{ }^{p} \longrightarrow[-\infty,+\infty[\right.\right.$, then

$$
\chi\left(u_{1}, \ldots, u_{p}\right) \in \operatorname{Sh}(\Omega) .
$$

In particular $u_{1}+\cdots+u_{p}, \max \left\{u_{1}, \ldots, u_{p}\right\}, \log \left(e^{u_{1}}+\cdots+e^{u_{p}}\right) \in \operatorname{Sh}(\Omega)$.
Proof. Every convex function is continuous, hence $\chi\left(u_{1}, \ldots, u_{p}\right)$ is upper semicontinuous. One can write

$$
\chi(t)=\sup _{i \in I} A_{i}(t)
$$

where $A_{i}(t)=a_{1} t_{1}+\cdots+a_{p} t_{p}+b$ is the family of affine functions that define supporting hyperplanes of the graph of $\chi$. As $\chi\left(t_{1}, \ldots, t_{p}\right)$ is non-decreasing in each $t_{j}$, we have $a_{j} \geqslant 0$, thus

$$
\sum_{1 \leqslant j \leqslant p} a_{j} u_{j}(x)+b \leqslant \mu_{B}\left(\sum a_{j} u_{j}+b ; x, r\right) \leqslant \mu_{B}\left(\chi\left(u_{1}, \ldots, u_{p}\right) ; x, r\right)
$$

for every ball $\bar{B}(x, r) \subset \Omega$. If one takes the supremum of this inequality over all the $A_{i}$ 's, it follows that $\chi\left(u_{1}, \ldots, u_{p}\right)$ satisfies the mean value inequality 4.12 c$)$. In the last example, the function $\chi\left(t_{1}, \ldots, t_{p}\right)=\log \left(e^{t_{1}}+\cdots+e^{t_{p}}\right)$ is convex because

$$
\sum_{1 \leqslant j, k \leqslant p} \frac{\partial^{2} \chi}{\partial t_{j} \partial t_{k}} \xi_{j} \xi_{k}=e^{-\chi} \sum \xi_{j}^{2} e^{t_{j}}-e^{-2 \chi}\left(\sum \xi_{j} e^{t_{j}}\right)^{2}
$$

and $\left(\sum \xi_{j} e^{t_{j}}\right)^{2} \leqslant\left(\sum \xi_{j}^{2} e^{t_{j}}\right) e^{\chi}$ by the Cauchy-Schwarz inequality.
(4.17) Theorem. If $\Omega$ is connected and $u \in \operatorname{Sh}(\Omega)$, then either $u \equiv-\infty$ or $u \in L_{\mathrm{loc}}^{1}(\Omega)$.

Proof. Note that a subharmonic function is always locally bounded above. Let $W$ be the set of points $x \in \Omega$ such that $u$ is integrable in a neighborhood of $x$. Then $W$ is open by definition and $u>-\infty$ almost everywhere on $W$. If $x \in \bar{W}$, one can choose $a \in W$ such that $|a-x|<r=\frac{1}{2} d(x, \complement \Omega)$ and $u(a)>-\infty$. Then $B(a, r)$ is a neighborhood of $x, \bar{B}(a, r) \subset \Omega$ and $\mu_{B}(u ; a, r) \geqslant u(a)>-\infty$. Therefore $x \in W, W$ is also closed. We must have $W=\Omega$ or $W=\emptyset$; in the last case $u \equiv-\infty$ by the mean value inequality.
(4.18) Theorem. Let $u \in \operatorname{Sh}(\Omega)$ be such that $u \not \equiv-\infty$ on each connected component of $\Omega$. Then
a) $r \longmapsto \mu_{S}(u ; a, r), r \longmapsto \mu_{B}(u ; a, r)$ are non decreasing functions in the interval $] 0, d(a, \complement \Omega)\left[\right.$, and $\mu_{B}(u ; a, r) \leqslant \mu_{S}(u ; a, r)$.
b) For any family $\left(\rho_{\varepsilon}\right)$ of smoothing kernels, $u \star \rho_{\varepsilon} \in \operatorname{Sh}\left(\Omega_{\varepsilon}\right) \cap \mathscr{C}^{\infty}\left(\Omega_{\varepsilon}, \mathbb{R}\right)$, the family $\left(u \star \rho_{\varepsilon}\right)$ is non decreasing in $\varepsilon$ and $\lim _{\varepsilon \rightarrow 0} u \star \rho_{\varepsilon}=u$.

Proof. We first verify statements a) and b) when $u \in C^{2}(\Omega, \mathbb{R})$. Then $\Delta u \geqslant 0$ and $\mu_{S}(u ; a, r)$ is non decreasing in virtue of (4.10). By (4.9), we find that $\mu_{B}(u ; a, r)$ is also non decreasing and that $\mu_{B}(u ; a, r) \leqslant \mu_{S}(u ; a, r)$. Furthermore, Formula (4.11) shows that $\varepsilon \longmapsto u \star \rho_{\varepsilon}(a)$ is non decreasing (provided that $\rho_{\varepsilon}$ is radially symmetric).

In the general case, we first observe that property 4.12 c ) is equivalent to the inequality

$$
u \leqslant u \star \mu_{r} \quad \text { on } \quad \Omega_{r}, \quad \forall r>0,
$$

where $\mu_{r}$ is the probability measure of uniform density on $B(0, r)$. This inequality implies $u \star \rho_{\varepsilon} \leqslant u \star \rho_{\varepsilon} \star \mu_{r}$ on $\left(\Omega_{r}\right)_{\varepsilon}=\Omega_{r+\varepsilon}$, thus $u \star \rho_{\varepsilon} \in \mathscr{C}^{\infty}\left(\Omega_{\varepsilon}, \mathbb{R}\right)$ is subharmonic on $\Omega_{\varepsilon}$. It follows that $u \star \rho_{\varepsilon} \star \rho_{\eta}$ is non decreasing in $\eta$; by symmetry, it is also non decreasing in $\varepsilon$, and so is $u \star \rho_{\varepsilon}=\lim _{\eta \rightarrow 0} u \star \rho_{\varepsilon} \star \rho_{\eta}$. We have $u \star \rho_{\varepsilon} \geqslant u$ by (4.19) and $\lim \sup _{\varepsilon \rightarrow 0} u \star \rho_{\varepsilon} \leqslant u$ by the upper semicontinuity. Hence $\lim _{\varepsilon \rightarrow 0} u \star \rho_{\varepsilon}=u$. Property a) for $u$ follows now from its validity for $u \star \rho_{\varepsilon}$ and from the monotone convergence theorem.
(4.19) Corollary. If $u \in \operatorname{Sh}(\Omega)$ is such that $u \not \equiv-\infty$ on each connected component of $\Omega$, then $\Delta u$ computed in the sense of distribution theory is a positive measure.

Indeed $\Delta\left(u \star \rho_{\varepsilon}\right) \geqslant 0$ as a function, and $\Delta\left(u \star \rho_{\varepsilon}\right)$ converges weakly to $\Delta u$ in $\mathscr{D}^{\prime}(\Omega)$. Corollary 4.19 has a converse, but the correct statement is slightly more involved than for the direct property:
(4.20) Theorem. If $v \in \mathscr{D}^{\prime}(\Omega)$ is such that $\Delta v$ is a positive measure, there exists a unique function $u \in \operatorname{Sh}(\Omega)$ locally integrable such that $v$ is the distribution associated to $u$.

We must point out that $u$ need not coincide everywhere with $v$, even when $v$ is a locally integrable upper semicontinuous function: for example, if $v$ is the characteristic function of a compact subset $K \subset \Omega$ of measure 0 , the subharmonic representant of $v$ is $u=0$.

Proof. Set $v_{\varepsilon}=v \star \rho_{\varepsilon} \in \mathscr{C}^{\infty}\left(\Omega_{\varepsilon}, \mathbb{R}\right)$. Then $\Delta v_{\varepsilon}=(\Delta v) \star \rho_{\varepsilon} \geqslant 0$, thus $v_{\varepsilon} \in \operatorname{Sh}\left(\Omega_{\varepsilon}\right)$. Arguments similar to those in the proof of Th. 4.18 show that $\left(v_{\varepsilon}\right)$ is non decreasing in $\varepsilon$. Then $u:=\lim _{\varepsilon \rightarrow 0} v_{\varepsilon} \in \operatorname{Sh}(\Omega)$ by Th. 4.15. Since $v_{\varepsilon}$ converges weakly to $v$, the monotone convergence theorem shows that

$$
\langle v, f\rangle=\lim _{\varepsilon \rightarrow 0} \int_{\Omega} v_{\varepsilon} f d \lambda=\int_{\Omega} u f d \lambda, \quad \forall f \in \mathscr{D}(\Omega), \quad f \geqslant 0
$$

which concludes the existence part. The uniqueness of $u$ is clear from the fact that $u$ must satisfy $u=\lim u \star \rho_{\varepsilon}=\lim v \star \rho_{\varepsilon}$.

The most natural topology on the space $\operatorname{Sh}(\Omega)$ of subharmonic functions is the topology induced by the vector space topology of $L_{\mathrm{loc}}^{1}(\Omega)$ (Fréchet topology of convergence in $L^{1}$ norm on every compact subset of $\Omega$ ).
(4.21) Proposition. The convex cone $\operatorname{Sh}(\Omega) \cap L_{\mathrm{loc}}^{1}(\Omega)$ is closed in $L_{\mathrm{loc}}^{1}(\Omega)$, and it has the property that every bounded subset is relatively compact.

Proof. Let $\left(u_{j}\right)$ be a sequence in $\operatorname{Sh}(\Omega) \cap L_{\text {loc }}^{1}(\Omega)$. If $u_{j} \rightarrow u$ in $L_{\mathrm{loc}}^{1}(\Omega)$ then $\Delta u_{j} \rightarrow \Delta u$ in the weak topology of distributions, hence $\Delta u \geqslant 0$ and $u$ can be represented by a subharmonic function thanks to Th. 4.20. Now, suppose that $\left\|u_{j}\right\|_{L^{1}(K)}$ is uniformly bounded for every compact subset $K$ of $\Omega$. Let $\mu_{j}=\Delta u_{j} \geqslant 0$. If $\psi \in \mathscr{D}(\Omega)$ is a test function equal to 1 on a neighborhood $\omega$ of $K$ and such that $0 \leqslant \psi \leqslant 1$ on $\Omega$, we find

$$
\mu_{j}(K) \leqslant \int_{\Omega} \psi \Delta u_{j} d \lambda=\int_{\Omega} \Delta \psi u_{j} d \lambda \leqslant C\left\|u_{j}\right\|_{L^{1}\left(K^{\prime}\right)}
$$

where $K^{\prime}=\operatorname{Supp} \psi$, hence the sequence of measures $\left(\mu_{j}\right)$ is uniformly bounded in mass on every compact subset of $\Omega$. By weak compactness, there is a subsequence ( $\mu_{j_{\nu}}$ ) which converges weakly to a positive measure $\mu$ on $\Omega$. We claim that $f \star\left(\psi \mu_{j_{\nu}}\right)$ converges to $f \star(\psi \mu)$ in $L_{\text {loc }}^{1}\left(\mathbb{R}^{m}\right)$ for every function $f \in L_{\text {loc }}^{1}\left(\mathbb{R}^{m}\right)$. In fact, this is clear if $f \in \mathscr{C}^{\infty}\left(\mathbb{R}^{m}\right)$, and in general we use an approximation of $f$ by a smooth function $g$ together with the estimate

$$
\left\|(f-g) \star\left(\psi \mu_{j_{\nu}}\right)\right\|_{L^{1}(A)} \leqslant\|(f-g)\|_{L^{1}\left(A+K^{\prime}\right)} \mu_{j_{\nu}}\left(K^{\prime}\right), \quad \forall A \subset \subset \mathbb{R}^{m}
$$

to get the conclusion. We apply this when $f=N$ is the Newton kernel. Then $h_{j}=$ $u_{j}-N \star\left(\psi \mu_{j}\right)$ is harmonic on $\omega$ and bounded in $L^{1}(\omega)$. As $h_{j}=h_{j} \star \rho_{\varepsilon}$ for any smoothing kernel $\rho_{\varepsilon}$, we see that all derivatives $D^{\alpha} h_{j}=h_{j} \star\left(D^{\alpha} \rho_{\varepsilon}\right)$ are in fact uniformly locally bounded in $\omega$. Hence, after extracting a new subsequence, we may suppose that
$h_{j_{\nu}}$ converges uniformly to a limit $h$ on $\omega$. Then $u_{j_{\nu}}=h_{j_{\nu}}+N \star\left(\psi \mu_{j_{\nu}}\right)$ converges to $u=h+N \star(\psi \mu)$ in $L_{\mathrm{loc}}^{1}(\omega)$, as desired.

We conclude this subsection by stating a generalized version of the Green-Riesz formula.
(4.22) Proposition. Let $u \in \operatorname{Sh}(\Omega) \cap L_{\mathrm{loc}}^{1}(\Omega)$ and $\bar{B}(0, r) \subset \Omega$.
a) The Green-Riesz formula still holds true for such an $u$, namely, for every $x \in B(0, r)$

$$
u(x)=\int_{B(0, r)} \Delta u(y) G_{r}(x, y) d \lambda(y)+\int_{S(0, r)} u(y) P_{r}(x, y) d \sigma(y)
$$

b) (Harnack inequality)

If $u \geqslant 0$ on $\bar{B}(0, r)$, then for all $x \in B(0, r)$

$$
0 \leqslant u(x) \leqslant \int_{S(0, r)} u(y) P_{r}(x, y) d \sigma(y) \leqslant \frac{r^{m-2}(r+|x|)}{(r-|x|)^{m-1}} \mu_{S}(u ; 0, r)
$$

If $u \leqslant 0$ on $\bar{B}(0, r)$, then for all $x \in B(0, r)$

$$
u(x) \leqslant \int_{S(0, r)} u(y) P_{r}(x, y) d \sigma(y) \leqslant \frac{r^{m-2}(r-|x|)}{(r+|x|)^{m-1}} \mu_{S}(u ; 0, r) \leqslant 0
$$

Proof. We know that a) holds true if $u$ is of class $C^{2}$. In general, we replace $u$ by $u \star \rho_{\varepsilon}$ and take the limit. We only have to check that

$$
\int_{B(0, r)} \mu \star \rho_{\varepsilon}(y) G_{r}(x, y) d \lambda(y)=\lim _{\varepsilon \rightarrow 0} \int_{B(0, r)} \mu(y) G_{r}(x, y) d \lambda(y)
$$

for the positive measure $\mu=\Delta u$. Let us denote by $\widetilde{G}_{x}(y)$ the function such that

$$
\widetilde{G}_{x}(y)= \begin{cases}G_{r}(x, y) & \text { if } x \in B(0, r) \\ 0 & \text { if } x \notin B(0, r) .\end{cases}
$$

Then

$$
\begin{aligned}
\int_{B(0, r)} \mu \star \rho_{\varepsilon}(y) G_{r}(x, y) d \lambda(y) & =\int_{\mathbb{R}^{m}} \mu \star \rho_{\varepsilon}(y) \widetilde{G}_{x}(y) d \lambda(y) \\
& =\int_{\mathbb{R}^{m}} \mu(y) \widetilde{G}_{x} \star \rho_{\varepsilon}(y) d \lambda(y) .
\end{aligned}
$$

However $\widetilde{G}_{x}$ is continuous on $\mathbb{R}^{m} \backslash\{x\}$ and subharmonic in a neighborhood of $x$, hence $\widetilde{G}_{x} \star \rho_{\varepsilon}$ converges uniformly to $\widetilde{G}_{x}$ on every compact subset of $\mathbb{R}^{m} \backslash\{x\}$, and converges pointwise monotonically in a neighborhood of $x$. The desired equality follows by the monotone convergence theorem. Finally, b) is a consequence of a), for the integral involving $\Delta u$ is nonpositive and

$$
\frac{1}{\sigma_{m-1} r^{m-1}} \frac{r^{m-2}(r-|x|)}{(r+|x|)^{m-1}} \leqslant P_{r}(x, y) \leqslant \frac{1}{\sigma_{m-1} r^{m-1}} \frac{r^{m-2}(r+|x|)}{(r-|x|)^{m-1}}
$$

by (4.6) combined with the obvious inequality $(r-|x|)^{m} \leqslant|x-y|^{m} \leqslant(r+|x|)^{m}$.
$\S$ 4.C.3. Upper Envelopes and Choquet's Lemma. Let $\Omega \subset \mathbb{R}^{n}$ and let $\left(u_{\alpha}\right)_{\alpha \in I}$ be a family of upper semicontinuous functions $\Omega \longrightarrow\left[-\infty,+\infty\left[\right.\right.$. We assume that $\left(u_{\alpha}\right)$ is locally uniformly bounded above. Then the upper envelope

$$
u=\sup u_{\alpha}
$$

need not be upper semicontinuous, so we consider its upper semicontinuous regularization:

$$
u^{\star}(z)=\lim _{\varepsilon \rightarrow 0} \sup _{B(z, \varepsilon)} u \geqslant u(z) .
$$

It is easy to check that $u^{\star}$ is the smallest upper semicontinuous function which is $\geqslant u$. Our goal is to show that $u^{\star}$ can be computed with a countable subfamily of $\left(u_{\alpha}\right)$. Let $B\left(z_{j}, \varepsilon_{j}\right)$ be a countable basis of the topology of $\Omega$. For each $j$, let $\left(z_{j k}\right)$ be a sequence of points in $B\left(z_{j}, \varepsilon_{j}\right)$ such that

$$
\sup _{k} u\left(z_{j k}\right)=\sup _{B\left(z_{j}, \varepsilon_{j}\right)} u,
$$

and for each pair $(j, k)$, let $\alpha(j, k, l)$ be a sequence of indices $\alpha \in I$ such that $u\left(z_{j k}\right)=$ $\sup _{l} u_{\alpha(j, k, l)}\left(z_{j k}\right)$. Set

$$
v=\sup _{j, k, l} u_{\alpha(j, k, l)} .
$$

Then $v \leqslant u$ and $v^{\star} \leqslant u^{\star}$. On the other hand

$$
\sup _{B\left(z_{j}, \varepsilon_{j}\right)} v \geqslant \sup _{k} v\left(z_{j k}\right) \geqslant \sup _{k, l} u_{\alpha(j, k, l)}\left(z_{j k}\right)=\sup _{k} u\left(z_{j k}\right)=\sup _{B\left(z_{j}, \varepsilon_{j}\right)} u .
$$

As every ball $B(z, \varepsilon)$ is a union of balls $B\left(z_{j}, \varepsilon_{j}\right)$, we easily conclude that $v^{\star} \geqslant u^{\star}$, hence $v^{\star}=u^{\star}$. Therefore:
(4.23) Choquet's lemma. Every family $\left(u_{\alpha}\right)$ has a countable subfamily $\left(v_{j}\right)=\left(u_{\alpha(j)}\right)$ such that its upper envelope $v$ satisfies $v \leqslant u \leqslant u^{\star}=v^{\star}$.
(4.24) Proposition. If all $u_{\alpha}$ are subharmonic, the upper regularization $u^{\star}$ is subharmonic and equal almost everywhere to $u$.

Proof. By Choquet's lemma we may assume that $\left(u_{\alpha}\right)$ is countable. Then $u=\sup u_{\alpha}$ is a Borel function. As each $u_{\alpha}$ satisfies the mean value inequality on every ball $\bar{B}(z, r) \subset \Omega$, we get

$$
u(z)=\sup u_{\alpha}(z) \leqslant \sup \mu_{B}\left(u_{\alpha} ; z, r\right) \leqslant \mu_{B}(u ; z, r) .
$$

The right-hand side is a continuous function of $z$, so we infer

$$
u^{\star}(z) \leqslant \mu_{B}(u ; z, r) \leqslant \mu_{B}\left(u^{\star} ; z, r\right)
$$

and $u^{\star}$ is subharmonic. By the upper semicontinuity of $u^{\star}$ and the above inequality we find $u^{\star}(z)=\lim _{r \rightarrow 0} \mu_{B}(u ; z, r)$, thus $u^{\star}=u$ almost everywhere by Lebesgue's lemma.

## § 5. Plurisubharmonic Functions

## §5.A. Definition and Basic Properties

Plurisubharmonic functions have been introduced independently by [Lelong 1942] and [Oka 1942] for the study of holomorphic convexity. They are the complex counterparts of subharmonic functions.
(5.1) Definition. A function $u: \Omega \longrightarrow\left[-\infty,+\infty\left[\right.\right.$ defined on an open subset $\Omega \subset \mathbb{C}^{n}$ is said to be plurisubharmonic if
a) $u$ is upper semicontinuous ;
b) for every complex line $L \subset \mathbb{C}^{n}$, $u_{\mid \Omega \cap L}$ is subharmonic on $\Omega \cap L$.

The set of plurisubharmonic functions on $\Omega$ is denoted by $\operatorname{Psh}(\Omega)$.
An equivalent way of stating property b) is: for all $a \in \Omega, \xi \in \mathbb{C}^{n},|\xi|<d(a, C \Omega)$, then

$$
\begin{equation*}
u(a) \leqslant \frac{1}{2 \pi} \int_{0}^{2 \pi} u\left(a+e^{\mathrm{i} \theta} \xi\right) d \theta \tag{5.2}
\end{equation*}
$$

An integration of (5.2) over $\xi \in S(0, r)$ yields $u(a) \leqslant \mu_{S}(u ; a, r)$, therefore

$$
\begin{equation*}
\operatorname{Psh}(\Omega) \subset \operatorname{Sh}(\Omega) \tag{5.3}
\end{equation*}
$$

The following results have already been proved for subharmonic functions and are easy to extend to the case of plurisubharmonic functions:
(5.4) Theorem. For any decreasing sequence of plurisubharmonic functions $u_{k} \in$ $\operatorname{Psh}(\Omega)$, the limit $u=\lim u_{k}$ is plurisubharmonic on $\Omega$.
(5.5) Theorem. Let $u \in \operatorname{Psh}(\Omega)$ be such that $u \not \equiv-\infty$ on every connected component of $\Omega$. If $\left(\rho_{\varepsilon}\right)$ is a family of smoothing kernels, then $u \star \rho_{\varepsilon}$ is $\mathscr{C}^{\infty}$ and plurisubharmonic on $\Omega_{\varepsilon}$, the family $\left(u \star \rho_{\varepsilon}\right)$ is non decreasing in $\varepsilon$ and $\lim _{\varepsilon \rightarrow 0} u \star \rho_{\varepsilon}=u$.
(5.6) Theorem. Let $u_{1}, \ldots, u_{p} \in \operatorname{Psh}(\Omega)$ and $\chi: \mathbb{R}^{p} \longrightarrow \mathbb{R}$ be a convex function such that $\chi\left(t_{1}, \ldots, t_{p}\right)$ is non decreasing in each $t_{j}$. Then $\chi\left(u_{1}, \ldots, u_{p}\right)$ is plurisubharmonic on $\Omega$. In particular $u_{1}+\cdots+u_{p}$, max $\left\{u_{1}, \ldots, u_{p}\right\}, \log \left(e^{u_{1}}+\cdots+e^{u_{p}}\right)$ are plurisubharmonic on $\Omega$.
(5.7) Theorem. Let $\left\{u_{\alpha}\right\} \subset \operatorname{Psh}(\Omega)$ be locally uniformly bounded from above and $u=\sup u_{\alpha}$. Then the regularized upper envelope $u^{\star}$ is plurisubharmonic and is equal to $u$ almost everywhere.

Proof. By Choquet's lemma, we may assume that $\left(u_{\alpha}\right)$ is countable. Then $u$ is a Borel function which clearly satisfies (5.2), and thus $u \star \rho_{\varepsilon}$ also satisfies (5.2). Hence $u \star$ $\rho_{\varepsilon}$ is plurisubharmonic. By Proposition 4.24, $u^{\star}=u$ almost everywhere and $u^{\star}$ is subharmonic, so

$$
u^{\star}=\lim u^{\star} \star \rho_{\varepsilon}=\lim u \star \rho_{\varepsilon}
$$

is plurisubharmonic.
If $u \in C^{2}(\Omega, \mathbb{R})$, the subharmonicity of restrictions of $u$ to complex lines, $\mathbb{C} \ni w \longmapsto$ $u(a+w \xi), a \in \Omega, \xi \in \mathbb{C}^{n}$, is equivalent to

$$
\frac{\partial^{2}}{\partial w \partial \bar{w}} u(a+w \xi)=\sum_{1 \leqslant j, k \leqslant n} \frac{\partial^{2} u}{\partial z_{j} \partial \bar{z}_{k}}(a+w \xi) \xi_{j} \bar{\xi}_{k} \geqslant 0 .
$$

Therefore, $u$ is plurisubharmonic on $\Omega$ if and only if $\sum \partial^{2} u / \partial z_{j} \partial \bar{z}_{k}(a) \xi_{j} \bar{\xi}_{k}$ is a semipositive hermitian form at every point $a \in \Omega$. This equivalence is still true for arbitrary plurisubharmonic functions, under the following form:
(5.8) Theorem. If $u \in \operatorname{Psh}(\Omega), u \not \equiv-\infty$ on every connected component of $\Omega$, then for all $\xi \in \mathbb{C}^{n}$

$$
H u(\xi):=\sum_{1 \leqslant j, k \leqslant n} \frac{\partial^{2} u}{\partial z_{j} \partial \bar{z}_{k}} \xi_{j} \bar{\xi}_{k} \in \mathscr{D}^{\prime}(\Omega)
$$

is a positive measure. Conversely, if $v \in \mathscr{D}^{\prime}(\Omega)$ is such that $H v(\xi)$ is a positive measure for every $\xi \in \mathbb{C}^{n}$, there exists a unique function $u \in \operatorname{Psh}(\Omega)$ locally integrable on $\Omega$ such that $v$ is the distribution associated to $u$.

Proof. If $u \in \operatorname{Psh}(\Omega)$, then $H u(\xi)=$ weak $\lim H\left(u \star \rho_{\varepsilon}\right)(\xi) \geqslant 0$. Conversely, $H v \geqslant 0$ implies $H\left(v \star \rho_{\varepsilon}\right)=(H v) \star \rho_{\varepsilon} \geqslant 0$, thus $v \star \rho_{\varepsilon} \in \operatorname{Psh}(\Omega)$, and also $\Delta v \geqslant 0$, hence $\left(v \star \rho_{\varepsilon}\right)$ is non decreasing in $\varepsilon$ and $u=\lim _{\varepsilon \rightarrow 0} v \star \rho_{\varepsilon} \in \operatorname{Psh}(\Omega)$ by Th. 5.4.
(5.9) Proposition. The convex cone $\operatorname{Psh}(\Omega) \cap L_{\mathrm{loc}}^{1}(\Omega)$ is closed in $L_{\mathrm{loc}}^{1}(\Omega)$, and it has the property that every bounded subset is relatively compact.

## § 5.B. Relations with Holomorphic Functions

In order to get a better geometric insight, we assume more generally that $u$ is a $C^{2}$ function on a complex $n$-dimensional manifold $X$. The complex Hessian of $u$ at a point $a \in X$ is the hermitian form on $T_{X}$ defined by

$$
\begin{equation*}
H u_{a}=\sum_{1 \leqslant j, k \leqslant n} \frac{\partial^{2} u}{\partial z_{j} \partial \bar{z}_{k}}(a) d z_{j} \otimes d \bar{z}_{k} \tag{5.10}
\end{equation*}
$$

If $F: X \longrightarrow Y$ is a holomorphic mapping and if $v \in C^{2}(Y, \mathbb{R})$, we have $d^{\prime} d^{\prime \prime}(v \circ F)=$ $F^{\star} d^{\prime} d^{\prime \prime} v$. In equivalent notations, a direct calculation gives for all $\xi \in T_{X, a}$

$$
H(v \circ F)_{a}(\xi)=\sum_{j, k, l, m} \frac{\partial^{2} v(F(a))}{\partial z_{l} \partial \bar{z}_{m}} \frac{\partial F_{l}(a)}{\partial z_{j}} \xi_{j} \overline{\frac{\partial F_{m}(a)}{\partial z_{k}} \xi_{k}}=H v_{F(a)}\left(F^{\prime}(a) . \xi\right)
$$

In particular $H u_{a}$ does not depend on the choice of coordinates $\left(z_{1}, \ldots, z_{n}\right)$ on $X$, and $H v_{a} \geqslant 0$ on $Y$ implies $H(v \circ F)_{a} \geqslant 0$ on $X$. Therefore, the notion of plurisubharmonic function makes sense on any complex manifold.
(5.11) Theorem. If $F: X \longrightarrow Y$ is a holomorphic map and $v \in \operatorname{Psh}(Y)$, then $v \circ F \in$ $\operatorname{Psh}(X)$.

Proof. It is enough to prove the result when $X=\Omega_{1} \subset \mathbb{C}^{n}$ and $X=\Omega_{2} \subset \mathbb{C}^{p}$ are open subsets. The conclusion is already known when $v$ is of class $C^{2}$, and it can be extended to an arbitrary upper semicontinuous function $v$ by using Th. 5.4 and the fact that $v=\lim v \star \rho_{\varepsilon}$.
(5.12) Example. By (3.22) we see that $\log |z|$ is subharmonic on $\mathbb{C}$, thus $\log |f| \in$ $\operatorname{Psh}(X)$ for every holomorphic function $f \in \mathscr{O}(X)$. More generally

$$
\log \left(\left|f_{1}\right|^{\alpha_{1}}+\cdots+\left|f_{q}\right|^{\alpha_{q}}\right) \in \operatorname{Psh}(X)
$$

for every $f_{j} \in \mathscr{O}(X)$ and $\alpha_{j} \geqslant 0$ (apply Th. 5.6 with $u_{j}=\alpha_{j} \log \left|f_{j}\right|$ ).

## § 5.C. Convexity Properties

The close analogy of plurisubharmonicity with the concept of convexity strongly suggests that there are deeper connections between these notions. We describe here a few elementary facts illustrating this philosophy. Another interesting connection between plurisubharmonicity and convexity will be seen in § 7.B (Kiselman's minimum principle).
(5.13) Theorem. If $\Omega=\omega+\mathrm{i} \omega^{\prime}$ where $\omega$, $\omega^{\prime}$ are open subsets of $\mathbb{R}^{n}$, and if $u(z)$ is a plurisubharmonic function on $\Omega$ that depends only on $x=\operatorname{Re} z$, then $\omega \ni x \longmapsto u(x)$ is convex.

Proof. This is clear when $u \in C^{2}(\Omega, \mathbb{R})$, for $\partial^{2} u / \partial z_{j} \partial \bar{z}_{k}=\frac{1}{4} \partial^{2} u / \partial x_{j} \partial x_{k}$. In the general case, write $u=\lim u \star \rho_{\varepsilon}$ and observe that $u \star \rho_{\varepsilon}(z)$ depends only on $x$.
(5.14) Corollary. If $u$ is a plurisubharmonic function in the open polydisk $D(a, R)=$ $\prod D\left(a_{j}, R_{j}\right) \subset \mathbb{C}^{n}$, then

$$
\begin{aligned}
\mu\left(u ; r_{1}, \ldots, r_{n}\right) & =\frac{1}{(2 \pi)^{n}} \int_{0}^{2 \pi} u\left(a_{1}+r_{1} e^{\mathrm{i} \theta_{1}}, \ldots, a_{n}+r_{n} e^{\mathrm{i} \theta_{n}}\right) d \theta_{1} \ldots d \theta_{n} \\
m\left(u ; r_{1}, \ldots, r_{n}\right) & =\sup _{z \in D(a, r)} u\left(z_{1}, \ldots, z_{n}\right), \quad r_{j}<R_{j}
\end{aligned}
$$

are convex functions of $\left(\log r_{1}, \ldots, \log r_{n}\right)$ that are non decreasing in each variable.
Proof. That $\mu$ is non decreasing follows from the subharmonicity of $u$ along every coordinate axis. Now, it is easy to verify that the functions

$$
\begin{aligned}
\widetilde{\mu}\left(z_{1}, \ldots, z_{n}\right) & =\frac{1}{(2 \pi)^{n}} \int_{0}^{2 \pi} u\left(a_{1}+e^{z_{1}} e^{\mathrm{i} \theta_{1}}, \ldots, a_{n}+e^{z_{n}} e^{\mathrm{i} \theta_{n}}\right) d \theta_{1} \ldots d \theta_{n} \\
\widetilde{m}\left(z_{1}, \ldots, z_{n}\right) & =\sup _{\left|w_{j}\right| \leqslant 1} u\left(a_{1}+e^{z_{1}} w_{1}, \ldots, a_{n}+e^{z_{n}} w_{n}\right)
\end{aligned}
$$

are upper semicontinuous, satisfy the mean value inequality, and depend only on $\operatorname{Re} z_{j} \in$ $] 0, \log R_{j}$. Therefore $\widetilde{\mu}$ and $\widetilde{M}$ are convex. Cor. 5.14 follows from the equalities

$$
\begin{aligned}
\mu\left(u ; r_{1}, \ldots, r_{n}\right) & =\widetilde{\mu}\left(\log r_{1}, \ldots, \log r_{n}\right) \\
m\left(u ; r_{1}, \ldots, r_{n}\right) & =\widetilde{m}\left(\log r_{1}, \ldots, \log r_{n}\right) .
\end{aligned}
$$

## § 5.D. Pluriharmonic Functions

Pluriharmonic functions are the counterpart of harmonic functions in the case of functions of complex variables:
(5.15) Definition. $A$ function $u$ is said to be pluriharmonic if $u$ and $-u$ are plurisubharmonic.

A pluriharmonic function is harmonic (in particular smooth) in any $\mathbb{C}$-analytic coordinate system, and is characterized by the condition $H u=0$, i.e. $d^{\prime} d^{\prime \prime} u=0$ or

$$
\partial^{2} u / \partial z_{j} \partial \bar{z}_{k}=0 \quad \text { for all } j, k
$$

If $f \in \mathscr{O}(X)$, it follows that the functions $\operatorname{Re} f, \operatorname{Im} f$ are pluriharmonic. Conversely:
(5.16) Theorem. If the De Rham cohomology group $H_{\mathrm{DR}}^{1}(X, \mathbb{R})$ is zero, every pluriharmonic function $u$ on $X$ can be written $u=\operatorname{Re} f$ where $f$ is a holomorphic function on $X$.

Proof. By hypothesis $H_{\mathrm{DR}}^{1}(X, \mathbb{R})=0, u \in \mathscr{C}^{\infty}(X)$ and $d\left(d^{\prime} u\right)=d^{\prime \prime} d^{\prime} u=0$, hence there exists $g \in \mathscr{C}^{\infty}(X)$ such that $d g=d^{\prime} u$. Then $d g$ is of type $(1,0)$, i.e. $g \in \mathscr{O}(X)$ and

$$
d(u-2 \operatorname{Re} g)=d(u-g-\bar{g})=\left(d^{\prime} u-d g\right)+\left(d^{\prime \prime} u-d \bar{g}\right)=0 .
$$

Therefore $u=\operatorname{Re}(2 g+C)$, where $C$ is a locally constant function.

## § 5.E. Global Regularization of Plurisubharmonic Functions

We now study a very efficient regularization and patching procedure for continuous plurisubharmonic functions, essentially due to [Richberg 1968]. The main idea is contained in the following lemma:
(5.17) Lemma. Let $u_{\alpha} \in \operatorname{Psh}\left(\Omega_{\alpha}\right)$ where $\Omega_{\alpha} \subset \subset X$ is a locally finite open covering of $X$. Assume that for every index $\beta$

$$
\limsup _{\zeta \rightarrow z} u_{\beta}(\zeta)<\max _{\Omega_{\alpha} \ni z}\left\{u_{\alpha}(z)\right\}
$$

at all points $z \in \partial \Omega_{\beta}$. Then the function

$$
u(z)=\max _{\Omega_{\alpha} \ni z} u_{\alpha}(z)
$$

is plurisubharmonic on $X$.
Proof. Fix $z_{0} \in X$. Then the indices $\beta$ such that $z_{0} \in \partial \Omega_{\beta}$ or $z_{0} \notin \bar{\Omega}_{\beta}$ do not contribute to the maximum in a neighborhood of $z_{0}$. Hence there is a a finite set $I$ of indices $\alpha$ such that $\Omega_{\alpha} \ni z_{0}$ and a neighborhood $V \subset \bigcap_{\alpha \in I} \Omega_{\alpha}$ on which $u(z)=\max _{\alpha \in I} u_{\alpha}(z)$. Therefore $u$ is plurisubharmonic on $V$.

The above patching procedure produces functions which are in general only continuous. When smooth functions are needed, one has to use a regularized max function. Let
$\theta \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$ be a nonnegative function with support in $[-1,1]$ such that $\int_{\mathbb{R}} \theta(h) d h=1$ and $\int_{\mathbb{R}} h \theta(h) d h=0$.
(5.18) Lemma. For arbitrary $\left.\eta=\left(\eta_{1}, \ldots, \eta_{p}\right) \in\right] 0,+\infty\left[{ }^{p}\right.$, the function

$$
M_{\eta}\left(t_{1}, \ldots, t_{p}\right)=\int_{\mathbb{R}^{n}} \max \left\{t_{1}+h_{1}, \ldots, t_{p}+h_{p}\right\} \prod_{1 \leqslant j \leqslant n} \theta\left(h_{j} / \eta_{j}\right) d h_{1} \ldots d h_{p}
$$

possesses the following properties:
a) $M_{\eta}\left(t_{1}, \ldots, t_{p}\right)$ is non decreasing in all variables, smooth and convex on $\mathbb{R}^{n}$;
b) $\max \left\{t_{1}, \ldots, t_{p}\right\} \leqslant M_{\eta}\left(t_{1}, \ldots, t_{p}\right) \leqslant \max \left\{t_{1}+\eta_{1}, \ldots, t_{p}+\eta_{p}\right\}$;
c) $M_{\eta}\left(t_{1}, \ldots, t_{p}\right)=M_{\left(\eta_{1}, \ldots, \widehat{\eta_{j}}, \ldots, \eta_{p}\right)}\left(t_{1}, \ldots, \widehat{t_{j}}, \ldots, t_{p}\right)$ if $t_{j}+\eta_{j} \leqslant \max _{k \neq j}\left\{t_{k}-\eta_{k}\right\}$;
d) $M_{\eta}\left(t_{1}+a, \ldots, t_{p}+a\right)=M_{\eta}\left(t_{1}, \ldots, t_{p}\right)+a, \quad \forall a \in \mathbb{R}$;
e) if $u_{1}, \ldots, u_{p}$ are plurisubharmonic and satisfy $H\left(u_{j}\right)_{z}(\xi) \geqslant \gamma_{z}(\xi)$ where $z \mapsto \gamma_{z}$ is a continuous hermitian form on $T_{X}$, then $u=M_{\eta}\left(u_{1}, \ldots, u_{p}\right)$ is plurisubharmonic and satisfies $H u_{z}(\xi) \geqslant \gamma_{z}(\xi)$.

Proof. The change of variables $h_{j} \mapsto h_{j}-t_{j}$ shows that $M_{\eta}$ is smooth. All properties are immediate consequences of the definition, except perhaps e). That $M_{\eta}\left(u_{1}, \ldots, u_{p}\right)$ is plurisubharmonic follows from a) and Th. 5.6. Fix a point $z_{0}$ and $\varepsilon>0$. All functions $u_{j}^{\prime}(z)=u_{j}(z)-\gamma_{z_{0}}\left(z-z_{0}\right)+\varepsilon\left|z-z_{0}\right|^{2}$ are plurisubharmonic near $z_{0}$. It follows that

$$
M_{\eta}\left(u_{1}^{\prime}, \ldots, u_{p}^{\prime}\right)=u-\gamma_{z_{0}}\left(z-z_{0}\right)+\varepsilon\left|z-z_{0}\right|^{2}
$$

is also plurisubharmonic near $z_{0}$. Since $\varepsilon>0$ was arbitrary, e) follows.
(5.19) Corollary. Let $u_{\alpha} \in \mathscr{C}^{\infty}\left(\bar{\Omega}_{\alpha}\right) \cap \operatorname{Psh}\left(\Omega_{\alpha}\right)$ where $\Omega_{\alpha} \subset \subset X$ is a locally finite open covering of $X$. Assume that $u_{\beta}(z)<\max \left\{u_{\alpha}(z)\right\}$ at every point $z \in \partial \Omega_{\beta}$, when $\alpha$ runs over the indices such that $\Omega_{\alpha} \ni z$. Choose a family $\left(\eta_{\alpha}\right)$ of positive numbers so small that $u_{\beta}(z)+\eta_{\beta} \leqslant \max _{\Omega_{\alpha} \ni z}\left\{u_{\alpha}(z)-\eta_{\alpha}\right\}$ for all $\beta$ and $z \in \partial \Omega_{\beta}$. Then the function defined by

$$
\widetilde{u}(z)=M_{\left(\eta_{\alpha}\right)}\left(u_{\alpha}(z)\right) \quad \text { for } \alpha \text { such that } \Omega_{\alpha} \ni z
$$

is smooth and plurisubharmonic on $X$.
(5.20) Definition. A function $u \in \operatorname{Psh}(X)$ is said to be strictly plurisubharmonic if $u \in L_{\mathrm{loc}}^{1}(X)$ and if for every point $x_{0} \in X$ there exists a neighborhood $\Omega$ of $x_{0}$ and $c>0$ such that $u(z)-c|z|^{2}$ is plurisubharmonic on $\Omega$, i.e. $\sum\left(\partial^{2} u / \partial z_{j} \partial \bar{z}_{k}\right) \xi_{j} \bar{\xi}_{k} \geqslant c|\xi|^{2}($ as distributions on $\Omega$ ) for all $\xi \in \mathbb{C}^{n}$.
(5.21) Theorem ([Richberg 1968]). Let $u \in \operatorname{Psh}(X)$ be a continuous function which is strictly plurisubharmonic on an open subset $\Omega \subset X$, with $H u \geqslant \gamma$ for some continuous positive hermitian form $\gamma$ on $\Omega$. For any continuous function $\lambda \in C^{0}(\Omega), \lambda>0$, there exists a plurisubharmonic function $\widetilde{u}$ in $C^{0}(X) \cap \mathscr{C}^{\infty}(\Omega)$ such that $u \leqslant \widetilde{u} \leqslant u+\lambda$ on $\Omega$
and $\widetilde{u}=u$ on $X \backslash \Omega$, which is strictly plurisubharmonic on $\Omega$ and satisfies $H \widetilde{u} \geqslant(1-\lambda) \gamma$. In particular, $\widetilde{u}$ can be chosen strictly plurisubharmonic on $X$ if $u$ has the same property.

Proof. Let $\left(\Omega_{\alpha}\right)$ be a locally finite open covering of $\Omega$ by relatively compact open balls contained in coordinate patches of $X$. Choose concentric balls $\Omega_{\alpha}^{\prime \prime} \subset \Omega_{\alpha}^{\prime} \subset \Omega_{\alpha}$ of respective radii $r_{\alpha}^{\prime \prime}<r_{\alpha}^{\prime}<r_{\alpha}$ and center $z=0$ in the given coordinates $z=\left(z_{1}, \ldots, z_{n}\right)$ near $\bar{\Omega}_{\alpha}$, such that $\Omega_{\alpha}^{\prime \prime}$ still cover $\Omega$. We set

$$
u_{\alpha}(z)=u \star \rho_{\varepsilon_{\alpha}}(z)+\delta_{\alpha}\left(r_{\alpha}^{\prime 2}-|z|^{2}\right) \quad \text { on } \quad \bar{\Omega}_{\alpha} .
$$

For $\varepsilon_{\alpha}<\varepsilon_{\alpha, 0}$ and $\delta_{\alpha}<\delta_{\alpha, 0}$ small enough, we have $u_{\alpha} \leqslant u+\lambda / 2$ and $H u_{\alpha} \geqslant(1-\lambda) \gamma$ on $\bar{\Omega}_{\alpha}$. Set

$$
\eta_{\alpha}=\delta_{\alpha} \min \left\{r_{\alpha}^{\prime 2}-r_{\alpha}^{\prime \prime 2},\left(r_{\alpha}^{2}-r_{\alpha}^{\prime 2}\right) / 2\right\}
$$

Choose first $\delta_{\alpha}<\delta_{\alpha, 0}$ such that $\eta_{\alpha}<\min _{\bar{\Omega}_{\alpha}} \lambda / 2$, and then $\varepsilon_{\alpha}<\varepsilon_{\alpha, 0}$ so small that $u \leqslant u \star \rho_{\varepsilon_{\alpha}}<u+\eta_{\alpha}$ on $\bar{\Omega}_{\alpha}$. As $\delta_{\alpha}\left(r^{\prime 2}-|z|^{2}\right)$ is $\leqslant-2 \eta_{\alpha}$ on $\partial \Omega_{\alpha}$ and $>\eta_{\alpha}$ on $\bar{\Omega}_{\alpha}^{\prime \prime}$, we have $u_{\alpha}<u-\eta_{\alpha}$ on $\partial \Omega_{\alpha}$ and $u_{\alpha}>u+\eta_{\alpha}$ on $\bar{\Omega}_{\alpha}^{\prime \prime}$, so that the condition required in Corollary 5.19 is satisfied. We define

$$
\widetilde{u}= \begin{cases}u & \text { on } X \backslash \Omega, \\ M_{\left(\eta_{\alpha}\right)}\left(u_{\alpha}\right) & \text { on } \Omega .\end{cases}
$$

By construction, $\widetilde{u}$ is smooth on $\Omega$ and satisfies $u \leqslant \widetilde{u} \leqslant u+\lambda, H u \geqslant(1-\lambda) \gamma$ thanks to 5.18 (b,e). In order to see that $\widetilde{u}$ is plurisubharmonic on $X$, observe that $\widetilde{u}$ is the uniform limit of $\widetilde{u}_{\alpha}$ with

$$
\widetilde{u}_{\alpha}=\max \left\{u, M_{\left(\eta_{1} \ldots \eta_{\alpha}\right)}\left(u_{1} \ldots u_{\alpha}\right)\right\} \quad \text { on } \quad \bigcup_{1 \leqslant \beta \leqslant \alpha} \Omega_{\beta}
$$

and $\widetilde{u}_{\alpha}=u$ on the complement.

## § 5.F. Polar and Pluripolar Sets.

Polar and pluripolar sets are sets of $-\infty$ poles of subharmonic and plurisubharmonic functions. Although these functions possess a large amount of flexibility, pluripolar sets have some properties which remind their loose relationship with holomorphic functions.
(5.22) Definition. $A$ set $A \subset \Omega \subset \mathbb{R}^{m}$ (resp. $A \subset X, \operatorname{dim}_{\mathbb{C}} X=n$ ) is said to be polar (resp. pluripolar) if for every point $x \in \Omega$ there exist a connected neighborhood $W$ of $x$ and $u \in \operatorname{Sh}(W)($ resp. $u \in \operatorname{Psh}(W)), u \not \equiv-\infty$, such that $A \cap W \subset\{x \in W ; u(x)=-\infty\}$.

Theorem 4.17 implies that a polar or pluripolar set is of zero Lebesgue measure. Now, we prove a simple extension theorem.
(5.23) Theorem. Let $A \subset \Omega$ be a closed polar set and $v \in \operatorname{Sh}(\Omega \backslash A)$ such that $v$ is bounded above in a neighborhood of every point of $A$. Then $v$ has a unique extension $\widetilde{v} \in \operatorname{Sh}(\Omega)$.

Proof. The uniqueness is clear because $A$ has zero Lebesgue measure. On the other hand, every point of $A$ has a neighborhood $W$ such that

$$
A \cap W \subset\{x \in W ; u(x)=-\infty\}, \quad u \in \operatorname{Sh}(W), \quad u \not \equiv-\infty .
$$

After shrinking $W$ and subtracting a constant to $u$, we may assume $u \leqslant 0$. Then for every $\varepsilon>0$ the function $v_{\varepsilon}=v+\varepsilon u \in \operatorname{Sh}(W \backslash A)$ can be extended as an upper semicontinuous on $W$ by setting $v_{\varepsilon}=-\infty$ on $A \cap W$. Moreover, $v_{\varepsilon}$ satisfies the mean value inequality $v_{\varepsilon}(a) \leqslant \mu_{S}\left(v_{\varepsilon} ; a, r\right)$ if $a \in W \backslash A, r<d(a, A \cup C W)$, and also clearly if $a \in A, r<d(a, \mathrm{C} W)$. Therefore $v_{\varepsilon} \in \operatorname{Sh}(W)$ and $\widetilde{v}=\left(\sup v_{\varepsilon}\right)^{\star} \in \operatorname{Sh}(W)$. Clearly $\widetilde{v}$ coincides with $v$ on $W \backslash A$. A similar proof gives:
(5.24) Theorem. Let $A$ be a closed pluripolar set in a complex analytic manifold $X$. Then every function $v \in \operatorname{Psh}(X \backslash A)$ that is locally bounded above near $A$ extends uniquely into a function $\widetilde{v} \in \operatorname{Psh}(X)$.
(5.25) Corollary. Let $A \subset X$ be a closed pluripolar set. Every holomorphic function $f \in \mathbb{O}(X \backslash A)$ that is locally bounded near $A$ extends to a holomorphic function $\tilde{f} \in \mathbb{O}(X)$.

Proof. Apply Th. 5.24 to $\pm \operatorname{Re} f$ and $\pm \operatorname{Im} f$. It follows that $\operatorname{Re} f$ and $\operatorname{Im} f$ have pluriharmonic extensions to $X$, in particular $f$ extends to $\tilde{f} \in \mathscr{C}^{\infty}(X)$. By density of $X \backslash A$, $d^{\prime \prime} \widetilde{f}=0$ on $X$.
(5.26) Corollary. Let $A \subset \Omega$ (resp. $A \subset X$ ) be a closed (pluri)polar set. If $\Omega$ (resp. $X$ ) is connected, then $\Omega \backslash A$ (resp. $X \backslash A$ ) is connected.

Proof. If $\Omega \backslash A$ (resp. $X \backslash A$ ) is a disjoint union $\Omega_{1} \cup \Omega_{2}$ of non empty open subsets, the function defined by $f \equiv 0$ on $\Omega_{1}, f \equiv 1$ on $\Omega_{2}$ would have a harmonic (resp. holomorphic) extension through $A$, a contradiction.

## § 6. Domains of Holomorphy and Stein Manifolds

## $\S$ 6.A. Domains of Holomorphy in $\mathbb{C}^{n}$. Examples

Loosely speaking, a domain of holomorphy is an open subset $\Omega$ in $\mathbb{C}^{n}$ such that there is no part of $\partial \Omega$ across which all functions $f \in \mathscr{O}(\Omega)$ can be extended. More precisely:
(6.1) Definition. Let $\Omega \subset \mathbb{C}^{n}$ be an open subset. $\Omega$ is said to be a domain of holomorphy if for every connected open set $U \subset \mathbb{C}^{n}$ which meets $\partial \Omega$ and every connected component $V$ of $U \cap \Omega$ there exists $f \in \mathscr{O}(\Omega)$ such that $f_{\mid V}$ has no holomorphic extension to $U$.

Under the hypotheses made on $U$, we have $\emptyset \neq \partial V \cap U \subset \partial \Omega$. In order to show that $\Omega$ is a domain of holomorphy, it is thus sufficient to find for every $z_{0} \in \partial \Omega$ a function $f \in \mathscr{O}(\Omega)$ which is unbounded near $z_{0}$.
(6.2) Examples. Every open subset $\Omega \subset \mathbb{C}$ is a domain of holomorphy (for any $z_{0} \in \partial \Omega$, $f(z)=\left(z-z_{0}\right)^{-1}$ cannot be extended at $\left.z_{0}\right)$. In $\mathbb{C}^{n}$, every convex open subset is a domain of holomorphy: if $\operatorname{Re}\left\langle z-z_{0}, \xi_{0}\right\rangle=0$ is a supporting hyperplane of $\partial \Omega$ at $z_{0}$, the function $f(z)=\left(\left\langle z-z_{0}, \xi_{0}\right\rangle\right)^{-1}$ is holomorphic on $\Omega$ but cannot be extended at $z_{0}$.
(6.3) Hartogs figure. Assume that $n \geqslant 2$. Let $\omega \subset \mathbb{C}^{n-1}$ be a connected open set and $\omega^{\prime} \subsetneq \omega$ an open subset. Consider the open sets in $\mathbb{C}^{n}$ :

$$
\begin{array}{ll}
\Omega=((D(R) \backslash \bar{D}(r)) \times \omega) \cup\left(D(R) \times \omega^{\prime}\right) & \text { (Hartogs figure) } \\
\widetilde{\Omega}=D(R) \times \omega & \text { (filled Hartogs figure) }
\end{array}
$$

where $0 \leqslant r<R$ and $D(r) \subset \mathbb{C}$ denotes the open disk of center 0 and radius $r$ in $\mathbb{C}$.


Fig. I-3 Hartogs figure

Then every function $f \in \mathscr{O}(\Omega)$ can be extended to $\widetilde{\Omega}=\omega \times D(R)$ by means of the Cauchy formula:

$$
\widetilde{f}\left(z_{1}, z^{\prime}\right)=\frac{1}{2 \pi \mathrm{i}} \int_{\left|\zeta_{1}\right|=\rho} \frac{f\left(\zeta_{1}, z^{\prime}\right)}{\zeta_{1}-z_{1}} d \zeta_{1}, \quad z \in \widetilde{\Omega}, \quad \max \left\{\left|z_{1}\right|, r\right\}<\rho<R
$$

In fact $\tilde{f} \in \mathscr{O}(D(R) \times \omega)$ and $\tilde{f}=f$ on $D(R) \times \omega^{\prime}$, so we must have $\tilde{f}=f$ on $\Omega$ since $\Omega$ is connected. It follows that $\Omega$ is not a domain of holomorphy. Let us quote two interesting consequences of this example.
(6.4) Corollary (Riemann's extension theorem). Let $X$ be a complex analytic manifold, and $S$ a closed submanifold of codimension $\geqslant 2$. Then every $f \in \mathscr{O}(X \backslash S)$ extends holomorphically to $X$.

Proof. This is a local result. We may choose coordinates $\left(z_{1}, \ldots, z_{n}\right)$ and a polydisk $D(R)^{n}$ in the corresponding chart such that $S \cap D(R)^{n}$ is given by equations $z_{1}=\ldots=$ $z_{p}=0, p=\operatorname{codim} S \geqslant 2$. Then, denoting $\omega=D(R)^{n-1}$ and $\omega^{\prime}=\omega \backslash\left\{z_{2}=\ldots=z_{p}=\right.$ $0\}$, the complement $D(R)^{n} \backslash S$ can be written as the Hartogs figure

$$
D(R)^{n} \backslash S=((D(R) \backslash\{0\}) \times \omega) \cup\left(D(R) \times \omega^{\prime}\right)
$$

It follows that $f$ can be extended to $\widetilde{\Omega}=D(R)^{n}$.

## §6.B. Holomorphic Convexity and Pseudoconvexity

Let $X$ be a complex manifold. We first introduce the notion of holomorphic hull of a compact set $K \subset X$. This can be seen somehow as the complex analogue of the notion of (affine) convex hull for a compact set in a real vector space. It is shown that domains of holomorphy in $\mathbb{C}^{n}$ are characterized a property of holomorphic convexity. Finally, we prove that holomorphic convexity implies pseudoconvexity - a complex analogue of the geometric notion of convexity.
(6.5) Definition. Let $X$ be a complex manifold and let $K$ be a compact subset of $X$. Then the holomorphic hull of $K$ in $X$ is defined to be

$$
\widehat{K}=\widehat{K}_{\overparen{O}(X)}=\left\{z \in X ;|f(z)| \leqslant \sup _{K}|f|, \forall f \in \mathbb{O}(X)\right\} .
$$

## (6.6) Elementary properties.

a) $\widehat{K}$ is a closed subset of $X$ containing $K$. Moreover we have

$$
\sup _{\widehat{K}}|f|=\sup _{K}|f|, \quad \forall f \in \mathbb{O}(X)
$$

hence $\widehat{\widehat{K}}=\widehat{K}$.
b) If $h: X \rightarrow Y$ is a holomorphic map and $K \subset X$ is a compact set, then $h\left(\widehat{K}_{\mathscr{O}(X)}\right) \subset$ $\widehat{h(K})_{\overparen{O}(Y)}$. In particular, if $X \subset Y$, then $\widehat{K}_{\overparen{O}(X)} \subset \widehat{K}_{\mathscr{O}(Y)} \cap X$. This is immediate from the definition.
c) $\widehat{K}$ contains the union of $K$ with all relatively compact connected components of $X \backslash K$ (thus $\widehat{K}$ "fills the holes" of $K$ ). In fact, for every connected component $U$ of $X \backslash K$ we have $\partial U \subset \partial K$, hence if $\bar{U}$ is compact the maximum principle yields

$$
\sup _{\bar{U}}|f|=\sup _{\partial U}|f| \leqslant \sup _{K}|f|, \quad \text { for all } f \in \mathbb{O}(X)
$$

d) More generally, suppose that there is a holomorphic map $h: U \rightarrow X$ defined on a relatively compact open set $U$ in a complex manifold $S$, such that $h$ extends as a continuous map $h: \bar{U} \rightarrow X$ and $h(\partial U) \subset K$. Then $h(\bar{U}) \subset \widehat{K}$. Indeed, for $f \in \mathscr{O}(X)$, the maximum principle again yields

$$
\sup _{\bar{U}}|f \circ h|=\sup _{\partial U}|f \circ h| \leqslant \sup _{K}|f| .
$$

This is especially useful when $U$ is the unit disk in $\mathbb{C}$.
e) Suppose that $X=\Omega \subset \mathbb{C}^{n}$ is an open set. By taking $f(z)=\exp (A(z))$ where $A$ is an arbitrary affine function, we see that $\widehat{K}_{\sigma(\Omega)}$ is contained in the intersection of all affine half-spaces containing $K$. Hence $\widehat{K}_{\overparen{O}(\Omega)}$ is contained in the affine convex hull $\widehat{K}_{\text {aff }}$. As a consequence $\widehat{K}_{\overparen{O}(\Omega)}$ is always bounded and $\widehat{K}_{\left.\overparen{( } \mathbb{C}^{n}\right)}$ is a compact set. However, when $\Omega$ is arbitrary, $\widehat{K}_{\mathcal{O}(\Omega)}$ is not always compact; for example, in case $\Omega=\mathbb{C}^{n} \backslash\{0\}, n \geqslant 2$, then $\mathscr{O}(\Omega)=\mathscr{O}\left(\mathbb{C}^{n}\right)$ and the holomorphic hull of $K=S(0,1)$ is the non compact set $\widehat{K}=\bar{B}(0,1) \backslash\{0\}$.
(6.7) Definition. A complex manifold $X$ is said to be holomorphically convex if the holomorphic hull $\widehat{K}_{O(X)}$ of every compact set $K \subset X$ is compact.
(6.8) Remark. A complex manifold $X$ is holomorphically convex if and only if there is an exhausting sequence of holomorphically compact subsets $K_{\nu} \subset X$, i.e. compact sets such that

$$
X=\bigcup K_{\nu}, \quad \widehat{K}_{\nu}=K_{\nu}, \quad K_{\nu}^{\circ} \supset K_{\nu-1}
$$

Indeed, if $X$ is holomorphically convex, we may define $K_{\nu}$ inductively by $K_{0}=\emptyset$ and $K_{\nu+1}=\left(K_{\nu}^{\prime} \cup L_{\nu}\right)_{\hat{0}(X)}^{\wedge}$, where $K_{\nu}^{\prime}$ is a neighborhood of $K_{\nu}$ and $L_{\nu}$ a sequence of compact sets of $X$ such that $X=\bigcup L_{\nu}$. The converse is obvious: if such a sequence $\left(K_{\nu}\right)$ exists, then every compact subset $K \subset X$ is contained in some $K_{\nu}$, hence $\widehat{K} \subset \widehat{K}_{\nu}=K_{\nu}$ is compact.

We now concentrate on domains of holomorphy in $\mathbb{C}^{n}$. We denote by $d$ and $B(z, r)$ the distance and the open balls associated to an arbitrary norm on $\mathbb{C}^{n}$, and we set for simplicity $B=B(0,1)$.
(6.9) Proposition. If $\Omega$ is a domain of holomorphy and $K \subset \Omega$ is a compact subset, then $d(\widehat{K}, \complement \Omega)=d(K, \complement \Omega)$ and $\widehat{K}$ is compact.

Proof. Let $f \in \mathscr{O}(\Omega)$. Given $r<d(K, \mathcal{C} \Omega)$, we denote by $M$ the supremum of $|f|$ on the compact subset $K+r \bar{B} \subset \Omega$. Then for every $z \in K$ and $\xi \in \bar{B}$, the function

$$
\begin{equation*}
\mathbb{C} \ni t \longmapsto f(z+t \xi)=\sum_{k=0}^{+\infty} \frac{1}{k!} D^{k} f(z)(\xi)^{k} t^{k} \tag{6.10}
\end{equation*}
$$

is analytic in the disk $|t|<r$ and bounded by $M$. The Cauchy inequalities imply

$$
\left|D^{k} f(z)(\xi)^{k}\right| \leqslant M k!r^{-k}, \quad \forall z \in K, \quad \forall \xi \in \bar{B} .
$$

As the left hand side is an analytic fuction of $z$ in $\Omega$, the inequality must also hold for $z \in \widehat{K}, \xi \in \bar{B}$. Every $f \in \mathscr{O}(\Omega)$ can thus be extended to any ball $B(z, r), z \in \widehat{K}$, by means of the power series (6.10). Hence $B(z, r)$ must be contained in $\Omega$, and this shows that $d(\widehat{K}, \complement \Omega) \geqslant r$. As $r<d(K, C \Omega)$ was arbitrary, we get $d(\widehat{K},\lceil\Omega) \geqslant d(K, C \Omega)$ and the converse inequality is clear, so $d(\widehat{K}, C \Omega)=d(K, C \Omega)$. As $\widehat{K}$ is bounded and closed in $\Omega$, this shows that $\widehat{K}$ is compact.
(6.11) Theorem. Let $\Omega$ be an open subset of $\mathbb{C}^{n}$. The following properties are equivalent:
a) $\Omega$ is a domain of holomorphy;
b) $\Omega$ is holomorphically convex;
c) For every countable subset $\left\{z_{j}\right\}_{j \in \mathbb{N}} \subset \Omega$ without accumulation points in $\Omega$ and every sequence of complex numbers $\left(a_{j}\right)$, there exists an interpolation function $F \in \mathscr{O}(\Omega)$ such that $F\left(z_{j}\right)=a_{j}$.
d) There exists a function $F \in \mathscr{O}(\Omega)$ which is unbounded on any neighborhood of any point of $\partial \Omega$.

Proof. d) $\Longrightarrow \mathrm{a}$ ) is obvious and a$) \Longrightarrow \mathrm{b}$ ) is a consequence of Prop. 6.9.
c) $\Longrightarrow \mathrm{d}$ ). If $\Omega=\mathbb{C}^{n}$ there is nothing to prove. Otherwise, select a dense sequence $\left(\zeta_{j}\right)$ in $\partial \Omega$ and take $z_{j} \in \Omega$ such that $d\left(z_{j}, \zeta_{j}\right)<2^{-j}$. Then the interpolation function $F \in \mathbb{O}(\Omega)$ such that $F\left(z_{j}\right)=j$ satisfies d).
b) $\Longrightarrow$ c). Let $K_{\nu} \subset \Omega$ be an exhausting sequence of holomorphically convex compact sets as in Remark 6.8. Let $\nu(j)$ be the unique index $\nu$ such that $z_{j} \in K_{\nu(j)+1} \backslash K_{\nu(j)}$. By the definition of a holomorphic hull, we can find a function $g_{j} \in \mathscr{O}(\Omega)$ such that

$$
\sup _{K_{\nu(j)}}\left|g_{j}\right|<\left|g_{j}\left(z_{j}\right)\right| .
$$

After multiplying $g_{j}$ by a constant, we may assume that $g_{j}\left(z_{j}\right)=1$. Let $P_{j} \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$ be a polynomial equal to 1 at $z_{j}$ and to 0 at $z_{0}, z_{1}, \ldots, z_{j-1}$. We set

$$
F=\sum_{j=0}^{+\infty} \lambda_{j} P_{j} g_{j}^{m_{j}}
$$

where $\lambda_{j} \in \mathbb{C}$ and $m_{j} \in \mathbb{N}$ are chosen inductively such that

$$
\begin{aligned}
& \lambda_{j}=a_{j}-\sum_{0 \leqslant k<j} \lambda_{k} P_{k}\left(z_{j}\right) g_{k}\left(z_{j}\right)^{m_{k}}, \\
& \left|\lambda_{j} P_{j} g_{j}^{m_{j}}\right| \leqslant 2^{-j} \quad \text { on } K_{\nu(j)}
\end{aligned}
$$

once $\lambda_{j}$ has been chosen, the second condition holds as soon as $m_{j}$ is large enough. Since $\left\{z_{j}\right\}$ has no accumulation point in $\Omega$, the sequence $\nu(j)$ tends to $+\infty$, hence the series converges uniformly on compact sets.

We now show that a holomorphically convex manifold must satisfy some more geometric convexity condition, known as pseudoconvexity, which is most easily described in terms of the existence of plurisubharmonic exhaustion functions.
(6.12) Definition. A function $\psi: X \longrightarrow[-\infty,+\infty[$ on a topological space $X$ is said to be an exhaustion if all sublevel sets $X_{c}:=\{z \in X ; \psi(z)<c\}, c \in \mathbb{R}$, are relatively compact. Equivalently, $\psi$ is an exhaustion if and only if $\psi$ tends to $+\infty$ relatively to the filter of complements $X \backslash K$ of compact subsets of $X$.

A function $\psi$ on an open set $\Omega \subset \mathbb{R}^{n}$ is thus an exhaustion if and only if $\psi(x) \rightarrow+\infty$ as $x \rightarrow \partial \Omega$ or $x \rightarrow \infty$. It is easy to check, cf. Exercise 8.8, that a connected open set $\Omega \subset \mathbb{R}^{n}$ is convex if and only if $\Omega$ has a locally convex exhaustion function. Since plurisubharmonic functions appear as the natural generalization of convex functions in complex analysis, we are led to the following definition.
(6.13) Definition. Let $X$ be a complex n-dimensional manifold. Then $X$ is said to be
a) weakly pseudoconvex if there exists a smooth plurisubharmonic exhaustion function $\psi \in \operatorname{Psh}(X) \cap \mathscr{C}^{\infty}(X) ;$
b) strongly pseudoconvex if there exists a smooth strictly plurisubharmonic exhaustion function $\psi \in \operatorname{Psh}(X) \cap \mathscr{C}^{\infty}(X)$, i.e. $H \psi$ is positive definite at every point.
(6.14) Theorem. Every holomorphically convex manifold $X$ is weakly pseudoconvex.

Proof. Let $\left(K_{\nu}\right)$ be an exhausting sequence of holomorphically convex compact sets as in Remark 6.8. For every point $a \in L_{\nu}:=K_{\nu+2} \backslash K_{\nu+1}^{\circ}$, one can select $g_{\nu, a} \in \mathbb{G}(\Omega)$ such
that $\sup _{K_{\nu}}\left|g_{\nu, a}\right|<1$ and $\left|g_{\nu, a}(a)\right|>1$. Then $\left|g_{\nu, a}(z)\right|>1$ in a neighborhood of $a ;$ by the Borel-Lebesgue lemma, one can find finitely many functions $\left(g_{\nu, a}\right)_{a \in I_{\nu}}$ such that

$$
\max _{a \in I_{\nu}}\left\{\left|g_{\nu, a}(z)\right|\right\}>1 \text { for } z \in L_{\nu}, \quad \max _{a \in I_{\nu}}\left\{\left|g_{\nu, a}(z)\right|\right\}<1 \text { for } z \in K_{\nu}
$$

For a sufficiently large exponent $p(\nu)$ we get

$$
\sum_{a \in I_{\nu}}\left|g_{\nu, a}\right|^{2 p(\nu)} \geqslant \nu \quad \text { on } L_{\nu}, \quad \sum_{a \in I_{\nu}}\left|g_{\nu, a}\right|^{2 p(\nu)} \leqslant 2^{-\nu} \text { on } K_{\nu} .
$$

It follows that the series

$$
\psi(z)=\sum_{\nu \in \mathbb{N}} \sum_{a \in I_{\nu}}\left|g_{\nu, a}(z)\right|^{2 p(\nu)}
$$

converges uniformly to a real analytic function $\psi \in \operatorname{Psh}(X)$ (see Exercise 8.11). By construction $\psi(z) \geqslant \nu$ for $z \in L_{\nu}$, hence $\psi$ is an exhaustion.
(6.15) Example. The converse to Theorem 6.14 does not hold. In fact let $X=$ $\mathbb{C}^{2} / \Gamma$ be the quotient of $\mathbb{C}^{2}$ by the free abelian group of rank 2 generated by the affine automorphisms

$$
g_{1}(z, w)=\left(z+1, e^{\mathrm{i} \theta_{1}} w\right), \quad g_{2}(z, w)=\left(z+\mathrm{i}, e^{\mathrm{i} \theta_{2}} w\right), \quad \theta_{1}, \theta_{2} \in \mathbb{R}
$$

Since $\Gamma$ acts properly discontinuously on $\mathbb{C}^{2}$, the quotient has a structure of a complex (non compact) 2-dimensional manifold. The function $w \mapsto|w|^{2}$ is $\Gamma$-invariant, hence it induces a function $\psi\left((z, w)^{\sim}\right)=|w|^{2}$ on $X$ which is in fact a plurisubharmonic exhaustion function. Therefore $X$ is weakly pseudoconvex. On the other hand, any holomorphic function $f \in \mathscr{O}(X)$ corresponds to a $\Gamma$-invariant holomorphic function $\widetilde{f}(z, w)$ on $\mathbb{C}^{2}$. Then $z \mapsto \widetilde{f}(z, w)$ is bounded for $w$ fixed, because $\widetilde{f}(z, w)$ lies in the image of the compact set $K \times \bar{D}(0,|w|), K=$ unit square in $\mathbb{C}$. By Liouville's theorem, $\widetilde{f}(z, w)$ does not depend on $z$. Hence functions $f \in \mathbb{O}(X)$ are in one-to-one correspondence with holomorphic functions $\widetilde{f}(w)$ on $\mathbb{C}$ such that $\widetilde{f}\left(e^{\mathrm{i} \theta_{j}} w\right)=\widetilde{f}(w)$. By looking at the Taylor expansion at the origin, we conclude that $\widetilde{f}$ must be a constant if $\theta_{1} \notin \mathbb{Q}$ or $\theta_{1} \notin \mathbb{Q}$ (if $\theta_{1}, \theta_{2} \in \mathbb{Q}$ and $m$ is the least common denominator of $\theta_{1}, \theta_{2}$, then $\tilde{f}$ is a power series of the form $\left.\sum \alpha_{k} w^{m k}\right)$. From this, it follows easily that $X$ is holomorphically convex if and only if $\theta_{1}, \theta_{2} \in \mathbb{Q}$.

## §6.C. Stein Manifolds

The class of holomorphically convex manifolds contains two types of manifolds of a rather different nature:

- domains of holomorphy $X=\Omega \subset \mathbb{C}^{n}$;
- compact complex manifolds.

In the first case we have a lot of holomorphic functions, in fact the functions in $\mathbb{O}(\Omega)$ separate any pair of points of $\Omega$. On the other hand, if $X$ is compact and connected, the sets $\operatorname{Psh}(X)$ and $\mathfrak{O}(X)$ consist of constant functions merely (by the maximum principle). It is therefore desirable to introduce a clear distinction between these two subclasses. For
this purpose, [Stein 1951] introduced the class of manifolds which are now called Stein manifolds.
(6.16) Definition. A complex manifold $X$ is said to be a Stein manifold if
a) $X$ is holomorphically convex;
b) $\mathfrak{O}(X)$ locally separates points in $X$, i.e. every point $x \in X$ has a neighborhood $V$ such that for any $y \in V \backslash\{x\}$ there exists $f \in \mathscr{O}(X)$ with $f(y) \neq f(x)$.

The second condition is automatic if $X=\Omega$ is an open subset of $\mathbb{C}^{n}$. Hence an open set $\Omega \subset \mathbb{C}^{n}$ is Stein if and only if $\Omega$ is a domain of holomorphy.
(6.17) Lemma. If a complex manifold $X$ satisfies the axiom ( 6.16 b ) of local separation, there exists a smooth nonnegative strictly plurisubharmonic function $u \in \operatorname{Psh}(X)$.

Proof. Fix $x_{0} \in X$. We first show that there exists a smooth nonnegative function $u_{0} \in \operatorname{Psh}(X)$ which is strictly plurisubharmonic on a neighborhood of $x_{0}$. Let $\left(z_{1}, \ldots, z_{n}\right)$ be local analytic coordinates centered at $x_{0}$, and if necessary, replace $z_{j}$ by $\lambda z_{j}$ so that the closed unit ball $\bar{B}=\left\{\sum\left|z_{j}\right|^{2} \leqslant 1\right\}$ is contained in the neighborhood $V \ni x_{0}$ on which $(6.16 \mathrm{~b})$ holds. Then, for every point $y \in \partial B$, there exists a holomorphic function $f \in$ $\mathfrak{O}(X)$ such that $f(y) \neq f\left(x_{0}\right)$. Replacing $f$ with $\lambda\left(f-f\left(x_{0}\right)\right)$, we can achieve $f\left(x_{0}\right)=0$ and $|f(y)|>1$. By compactness of $\partial B$, we find finitely many functions $f_{1}, \ldots, f_{N} \in \mathscr{G}(X)$ such that $v_{0}=\sum\left|f_{j}\right|^{2}$ satisfies $v_{0}\left(x_{0}\right)=0$, while $v_{0} \geqslant 1$ on $\partial B$. Now, we set

$$
u_{0}(z)= \begin{cases}v_{0}(z) & \text { on } X \backslash B, \\ M_{\varepsilon}\left\{v_{0}(z),\left(|z|^{2}+1\right) / 3\right\} & \text { on } B .\end{cases}
$$

where $M_{\varepsilon}$ are the regularized max functions defined in 5.18. Then $u_{0}$ is smooth and plurisubharmonic, coincides with $v_{0}$ near $\partial B$ and with $\left(|z|^{2}+1\right) / 3$ on a neighborhood of $x_{0}$. We can cover $X$ by countably many neighborhoods $\left(V_{j}\right)_{j \geqslant 1}$, for which we have a smooth plurisubharmonic functions $u_{j} \in \operatorname{Psh}(X)$ such that $u_{j}$ is strictly plurisubharmonic on $V_{j}$. Then select a sequence $\varepsilon_{j}>0$ converging to 0 so fast that $u=\sum \varepsilon_{j} u_{j} \in \mathscr{C}^{\infty}(X)$. The function $u$ is nonnegative and strictly plurisubharmonic everywhere on $X$.
(6.18) Theorem. Every Stein manifold is strongly pseudoconvex.

Proof. By Th. 6.14, there is a smooth exhaustion function $\psi \in \operatorname{Psh}(X)$. If $u \geqslant 0$ is strictly plurisubharmonic, then $\psi^{\prime}=\psi+u$ is a strictly plurisubharmonic exhaustion.

The converse problem to know whether every strongly pseudoconvex manifold is actually a Stein manifold is known as the Levi problem, and was raised by [Levi 1910] in the case of domains $\Omega \subset \mathbb{C}^{n}$. In that case, the problem has been solved in the affirmative independently by [Oka 1953], [Norguet 1954] and [Bremermann 1954]. The general solution of the Levi problem has been obtained by [Grauert 1958]. Our proof will rely on the theory of $L^{2}$ estimates for $d^{\prime \prime}$, which will be available only in Chapter VIII.


Fig. I-4 Hartogs figure with excrescence
(6.19) Remark. It will be shown later that Stein manifolds always have enough holomorphic functions to separate finitely many points, and one can even interpolate given values of a function and its derivatives of some fixed order at any discrete set of points. In particular, we might have replaced condition ( 6.16 b ) by the stronger requirement that $\mathscr{O}(X)$ separates any pair of points. On the other hand, there are examples of manifolds satisfying the local separation condition ( 6.16 b ), but not global separation. A simple example is obtained by attaching an excrescence inside a Hartogs figure, in such a way that the resulting map $\pi: X \rightarrow D=D(0,1)^{2}$ is not one-to-one (see Figure I-4 above); then $\mathfrak{O}(X)$ coincides with $\pi^{\star} \mathscr{O}(D)$.

## § 6.D. Heredity Properties

Holomorphic convexity and pseudoconvexity are preserved under quite a number of natural constructions. The main heredity properties can be summarized in the following Proposition.
(6.20) Proposition. Let $\mathscr{C}$ denote the class of holomorphically convex (resp. of Stein, or weakly pseudoconvex, strongly pseudoconvex manifolds).
a) If $X, Y \in \mathscr{C}$, then $X \times Y \in \mathscr{C}$.
b) If $X \in \mathscr{C}$ and $S$ is a closed complex submanifold of $X$, then $S \in \mathscr{C}$.
c) If $\left(S_{j}\right)_{1 \leqslant j \leqslant N}$ is a collection of (not necessarily closed) submanifolds of a complex manifold $X$ such that $S=\bigcap S_{j}$ is a submanifold of $X$, and if $S_{j} \in \mathscr{C}$ for all $j$, then $S \in \mathscr{C}$.
d) If $F: X \rightarrow Y$ is a holomorphic map and $S \subset X, S^{\prime} \subset Y$ are (not necessarily closed) submanifolds in the class $\mathscr{C}$, then $S \cap F^{-1}\left(S^{\prime}\right)$ is in $\mathscr{C}$, as long as it is a submanifold of $X$.
e) If $X$ is a weakly (resp. strongly) pseudoconvex manifold and $u$ is a smooth plurisubharmonic function on $X$, then the open set $\Omega=u^{-1}(]-\infty, c[$ is weakly (resp. strongly) pseudoconvex. In particular the sublevel sets

$$
X_{c}=\psi^{-1}(]-\infty, c[)
$$

of a (strictly) plurisubharmonic exhaustion function are weakly (resp. strongly) pseudoconvex.

Proof. All properties are more or less immediate to check, so we only give the main facts.
a) For $K \subset X, L \subset Y$ compact, we have $(K \times L)_{\hat{\theta}(X \times Y)}^{\wedge}=\widehat{K}_{\hat{( }(X)} \times \widehat{K}_{\hat{( }(Y)}$, and if $\varphi$, $\psi$ are plurisubharmonic exhaustions of $X, Y$, then $\varphi(x)+\psi(y)$ is a plurisubharmonic exhaustion of $X \times Y$.
b) For a compact set $K \subset S$, we have $\widehat{K}_{\overparen{\sigma(S)}} \subset \widehat{K}_{\mathscr{O}(X)} \cap S$, and if $\psi \in \operatorname{Psh}(X)$ is an exhaustion, then $\psi \upharpoonright S \in \operatorname{Psh}(S)$ is an exhaustion (since $S$ is closed).
c) $\bigcap S_{j}$ is a closed submanifold in $\prod S_{j}$ (equal to its intersection with the diagonal of $X^{N}$ ).
d) For a compact set $K \subset S \cap F^{-1}\left(S^{\prime}\right)$, we have

$$
\left.\widehat{K}_{\overparen{O}\left(S \cap F^{-1}\left(S^{\prime}\right)\right)} \subset \widehat{K}_{\overparen{O}(S)} \cap F^{-1}(\widehat{F(K})_{\overparen{O}\left(S^{\prime}\right)}\right)
$$

and if $\varphi, \psi$ are plurisubharmonic exhaustions of $S, S^{\prime}$, then $\varphi+\psi \circ F$ is a plurisubharmonic exhaustion of $S \cap F^{-1}\left(S^{\prime}\right)$.
e) $\varphi(z):=\psi(z)+1 /(c-u(z))$ is a (strictly) plurisubharmonic exhaustion function on $\Omega$.

## § 7. Pseudoconvex Open Sets in $\mathbb{C}^{n}$

## $\S$ 7.A. Geometric Characterizations of Pseudoconvex Open Sets

We first discuss some characterizations of pseudoconvex open sets in $\mathbb{C}^{n}$. We will need the following elementary criterion for plurisubharmonicity.
(7.1) Criterion. Let $v: \Omega \longrightarrow[-\infty,+\infty[$ be an upper semicontinuous function. Then $v$ is plurisubharmonic if and only if for every closed disk $\bar{\Delta}=z_{0}+\bar{D}(1) \eta \subset \Omega$ and every polynomial $P \in \mathbb{C}[t]$ such that $v\left(z_{0}+t \eta\right) \leqslant \operatorname{Re} P(t)$ for $|t|=1$, then $v\left(z_{0}\right) \leqslant \operatorname{Re} P(0)$.

Proof. The condition is necessary because $t \longmapsto v\left(z_{0}+t \eta\right)-\operatorname{Re} P(t)$ is subharmonic in a neighborhood of $\bar{D}(1)$, so it satisfies the maximum principle on $D(1)$ by Th. 4.14. Let us prove now the sufficiency. The upper semicontinuity of $v$ implies $v=\lim _{\nu \rightarrow+\infty} v_{\nu}$ on $\partial \Delta$ where $\left(v_{\nu}\right)$ is a strictly decreasing sequence of continuous functions on $\partial \Delta$. As trigonometric polynomials are dense in $C^{0}\left(S^{1}, \mathbb{R}\right)$, we may assume $v_{\nu}\left(z_{0}+e^{\mathrm{i} \theta} \eta\right)=\operatorname{Re} P_{\nu}\left(e^{\mathrm{i} \theta}\right)$, $P_{\nu} \in \mathbb{C}[t]$. Then $v\left(z_{0}+t \eta\right) \leqslant \operatorname{Re} P_{\nu}(t)$ for $|t|=1$, and the hypothesis implies

$$
v\left(z_{0}\right) \leqslant \operatorname{Re} P_{\nu}(0)=\frac{1}{2 \pi} \int_{0}^{2 \pi} \operatorname{Re} P_{\nu}\left(e^{\mathrm{i} \theta}\right) d \theta=\frac{1}{2 \pi} \int_{0}^{2 \pi} v_{\nu}\left(z_{0}+e^{\mathrm{i} \theta} \eta\right) d \theta
$$

Taking the limit when $\nu$ tends to $+\infty$ shows that $v$ satisfies the mean value inequality (5.2).

For any $z \in \Omega$ and $\xi \in \mathbb{C}^{n}$, we denote by

$$
\delta_{\Omega}(z, \xi)=\sup \{r>0 ; z+D(r) \xi \subset \Omega\}
$$

the distance from $z$ to $\partial \Omega$ in the complex direction $\xi$.
(7.2) Theorem. Let $\Omega \subset \mathbb{C}^{n}$ be an open subset. The following properties are equivalent:
a) $\Omega$ is strongly pseudoconvex (according to Def. 6.13 b );
b) $\Omega$ is weakly pseudoconvex;
c) $\Omega$ has a plurisubharmonic exhaustion function $\psi$.
d) $-\log \delta_{\Omega}(z, \xi)$ is plurisubharmonic on $\Omega \times \mathbb{C}^{n}$;
e) $-\log d(z, С \Omega)$ is plurisubharmonic on $\Omega$.

If one of these properties hold, $\Omega$ is said to be a pseudoconvex open set.
Proof. The implications a) $\Longrightarrow \mathrm{b}) \Longrightarrow \mathrm{c})$ are obvious. For the implication c) $\Longrightarrow$ d), we use Criterion 7.1. Consider a disk $\bar{\Delta}=\left(z_{0}, \xi_{0}\right)+\bar{D}(1)(\eta, \alpha)$ in $\Omega \times \mathbb{C}^{n}$ and a polynomial $P \in \mathbb{C}[t]$ such that

$$
-\log \delta_{\Omega}\left(z_{0}+t \eta, \xi_{0}+t \alpha\right) \leqslant \operatorname{Re} P(t) \quad \text { for } \quad|t|=1
$$

We have to verify that the inequality also holds when $|t|<1$. Consider the holomorphic mapping $h: \mathbb{C}^{2} \longrightarrow \mathbb{C}^{n}$ defined by

$$
h(t, w)=z_{0}+t \eta+w e^{-P(t)}\left(\xi_{0}+t \alpha\right) .
$$

By hypothesis

$$
\begin{aligned}
& h(\bar{D}(1) \times\{0\})=\operatorname{pr}_{1}(\bar{\Delta}) \subset \Omega \\
& h(\partial D(1) \times D(1)) \subset \Omega \quad\left(\text { since }\left|e^{-P}\right| \leqslant \delta_{\Omega} \quad \text { on } \partial \Delta\right),
\end{aligned}
$$

and the desired conclusion is that $h(\bar{D}(1) \times D(1)) \subset \Omega$. Let $J$ be the set of radii $r \geqslant 0$ such that $h(\bar{D}(1) \times \bar{D}(r)) \subset \Omega$. Then $J$ is an open interval $[0, R[, R>0$. If $R<1$, we get a contradiction as follows. Let $\psi \in \operatorname{Psh}(\Omega)$ be an exhaustion function and

$$
K=h(\partial D(1) \times \bar{D}(R)) \subset \subset \Omega, \quad c=\sup _{K} \psi .
$$

As $\psi \circ h$ is plurisubharmonic on a neighborhood of $\bar{D}(1) \times D(R)$, the maximum principle applied with respect to $t$ implies

$$
\psi \circ h(t, w) \leqslant c \quad \text { on } \bar{D}(1) \times D(R),
$$

hence $h(\bar{D}(1) \times D(R)) \subset \Omega_{c} \subset \subset \Omega$ and $h(\bar{D}(1) \times \bar{D}(R+\varepsilon)) \subset \Omega$ for some $\varepsilon>0$, a contradiction.
$\mathrm{d}) \Longrightarrow \mathrm{e})$. The function $-\log d(z, \complement \Omega)$ is continuous on $\Omega$ and satisfies the mean value inequality because

$$
-\log d(z, C \Omega)=\sup _{\xi \in \bar{B}}\left(-\log \delta_{\Omega}(z, \xi)\right) .
$$

e) $\Longrightarrow$ a). It is clear that

$$
u(z)=|z|^{2}+\max \left\{\log d(z, \complement \Omega)^{-1}, 0\right\}
$$

is a continuous strictly plurisubharmonic exhaustion function. Richberg's theorem 5.21 implies that there exists $\psi \in \mathscr{C}^{\infty}(\Omega)$ strictly plurisubharmonic such that $u \leqslant \psi \leqslant u+1$. Then $\psi$ is the required exhaustion function.

## (7.3) Proposition.

a) Let $\Omega \subset \mathbb{C}^{n}$ and $\Omega^{\prime} \subset \mathbb{C}^{p}$ be pseudoconvex. Then $\Omega \times \Omega^{\prime}$ is pseudoconvex. For every holomorphic map $F: \Omega \rightarrow \mathbb{C}^{p}$ the inverse image $F^{-1}\left(\Omega^{\prime}\right)$ is pseudoconvex.
b) If $\left(\Omega_{\alpha}\right)_{\alpha \in I}$ is a family of pseudoconvex open subsets of $\mathbb{C}^{n}$, the interior of the intersection $\Omega=\left(\bigcap_{\alpha \in I} \Omega_{\alpha}\right)^{\circ}$ is pseudoconvex.
c) If $\left(\Omega_{j}\right)_{j \in \mathbb{N}}$ is a non decreasing sequence of pseudoconvex open subsets of $\mathbb{C}^{n}$, then $\Omega=\bigcup_{j \in \mathbb{N}} \Omega_{j}$ is pseudoconvex.

Proof. a) Let $\varphi, \psi$ be smooth plurisubharmonic exhaustions of $\Omega, \Omega^{\prime}$. Then $(z, w) \longmapsto$ $\varphi(z)+\psi(w)$ is an exhaustion of $\Omega \times \Omega^{\prime}$ and $z \longmapsto \varphi(z)+\psi(F(z))$ is an exhaustion of $F^{-1}\left(\Omega^{\prime}\right)$.
b) We have $-\log d(z, \mathrm{C} \Omega)=\sup _{\alpha \in I}-\log d\left(z, \mathrm{C} \Omega_{\alpha}\right)$, so this function is plurisubharmonic.
c) The limit $-\log d(z, \complement \Omega)=\lim _{j \rightarrow+\infty}-\log d\left(z, \complement \Omega_{j}\right)$ is plurisubharmonic, hence $\Omega$ is pseudoconvex. This result cannot be generalized to strongly pseudoconvex manifolds: J.E. Fornaess in [Fornaess 1977] has constructed an increasing sequence of 2-dimensional Stein (even affine algebraic) manifolds $X_{\nu}$ whose union is not Stein; see Exercise 8.16.

## (7.4) Examples.

a) An analytic polyhedron in $\mathbb{C}^{n}$ is an open subset of the form

$$
P=\left\{z \in \mathbb{C}^{n} ;\left|f_{j}(z)\right|<1,1 \leqslant j \leqslant N\right\}
$$

where $\left(f_{j}\right)_{1 \leqslant j \leqslant N}$ is a family of analytic functions on $\mathbb{C}^{n}$. By 7.3 a), every analytic polyhedron is pseudoconvex.
b) Let $\omega \subset \mathbb{C}^{n-1}$ be pseudoconvex and let $u: \omega \longrightarrow[-\infty,+\infty[$ be an upper semicontinuous function. Then the Hartogs domain

$$
\Omega=\left\{\left(z_{1}, z^{\prime}\right) \in \mathbb{C} \times \omega ; \log \left|z_{1}\right|+u\left(z^{\prime}\right)<0\right\}
$$

is pseudoconvex if and and only if $u$ is plurisubharmonic. To see that the plurisubharmonicity of $u$ is necessary, observe that

$$
u\left(z^{\prime}\right)=-\log \delta_{\Omega}\left(\left(0, z^{\prime}\right),(1,0)\right)
$$

Conversely, assume that $u$ is plurisubharmonic and continuous. If $\psi$ is a plurisubharmonic exhaustion of $\omega$, then

$$
\psi\left(z^{\prime}\right)+|\log | z_{1}\left|+u\left(z^{\prime}\right)\right|^{-1}
$$

is an exhaustion of $\Omega$. This is no longer true if $u$ is not continuous, but in this case we may apply Property 7.3 c) to conclude that

$$
\Omega_{\varepsilon}=\left\{\left(z_{1}, z^{\prime}\right) ; d\left(z^{\prime},\lceil\omega)>\varepsilon, \log \left|z_{1}\right|+u \star \rho_{\varepsilon}\left(z^{\prime}\right)<0\right\}, \quad \Omega=\bigcup \Omega_{\varepsilon}\right.
$$

are pseudoconvex.
c) An open set $\Omega \subset \mathbb{C}^{n}$ is called a tube of base $\omega$ if $\Omega=\omega+\mathrm{i} \mathbb{R}^{n}$ for some open subset $\omega \subset \mathbb{R}^{n}$. Then of course $-\log d(z, \complement \Omega)=-\log (x, \complement \omega)$ depends only on the real part $x=\operatorname{Re} z$. By Th. 5.13 , this function is plurisubharmonic if and only if it is locally convex in $x$. Therefore $\Omega$ if pseudoconvex if and only if every connected component of $\omega$ is convex.
d) An open set $\Omega \subset \mathbb{C}^{n}$ is called a Reinhardt domain if $\left(e^{\mathrm{i} \theta_{1}} z_{1}, \ldots, e^{\mathrm{i} \theta_{n}} z_{n}\right)$ is in $\Omega$ for every $z=\left(z_{1}, \ldots, z_{n}\right) \in \Omega$ and $\theta_{1}, \ldots, \theta_{n} \in \mathbb{R}^{n}$. For such a domain, we consider the logarithmic indicatrix

$$
\omega^{\star}=\Omega^{\star} \cap \mathbb{R}^{n} \quad \text { with } \quad \Omega^{\star}=\left\{\zeta \in \mathbb{C}^{n} ;\left(e^{\zeta_{1}}, \ldots, e^{\zeta_{n}}\right) \in \Omega\right\} .
$$

It is clear that $\Omega^{\star}$ is a tube of base $\omega^{\star}$. Therefore every connected component of $\omega^{\star}$ must be convex if $\Omega$ is pseudoconvex. The converse is not true: $\Omega=\mathbb{C}^{n} \backslash\{0\}$ is not pseudoconvex for $n \geqslant 2$ although $\omega^{\star}=\mathbb{R}^{n}$ is convex. However, the Reinhardt open set

$$
\Omega^{\bullet}=\left\{\left(z_{1}, \ldots, z_{n}\right) \in(\mathbb{C} \backslash\{0\})^{n} ;\left(\log \left|z_{1}\right|, \ldots, \log \left|z_{n}\right|\right) \in \omega^{\star}\right\} \subset \Omega
$$

is easily seen to be pseudoconvex if $\omega^{\star}$ is convex: if $\chi$ is a convex exhaustion of $\omega^{\star}$, then $\psi(z)=\chi\left(\log \left|z_{1}\right|, \ldots, \log \left|z_{n}\right|\right)$ is a plurisubharmonic exhaustion of $\Omega^{\bullet}$. Similarly, if $\omega^{\star}$ is convex and such that $x \in \omega^{\star} \Longrightarrow y \in \omega^{\star}$ for $y_{j} \leqslant x_{j}$, we can take $\chi$ increasing in all variables and tending to $+\infty$ on $\partial \omega^{\star}$, hence the set

$$
\widetilde{\Omega}=\left\{\left(z_{1}, \ldots, z_{n}\right) \in \mathbb{C}^{n} ;\left|z_{j}\right| \leqslant e^{x_{j}} \text { for some } x \in \omega^{\star}\right\}
$$

is a pseudoconvex Reinhardt open set containing 0 .

## § 7.B. Kiselman's Minimum Principle

We already know that a maximum of plurisubharmonic functions is plurisubharmonic. However, if $v$ is a plurisubharmonic function on $X \times \mathbb{C}^{n}$, the partial minimum function on $X$ defined by $u(\zeta)=\inf _{z \in \Omega} v(\zeta, z)$ need not be plurisubharmonic. A simple counterexample in $\mathbb{C} \times \mathbb{C}$ is given by

$$
v(\zeta, z)=|z|^{2}+2 \operatorname{Re}(z \zeta)=|z+\bar{\zeta}|^{2}-|\zeta|^{2}, \quad u(\zeta)=-|\zeta|^{2} .
$$

It follows that the image $F(\Omega)$ of a pseudoconvex open set $\Omega$ by a holomorphic map $F$ need not be pseudoconvex. In fact, if

$$
\Omega=\left\{(t, \zeta, z) \in \mathbb{C}^{3} ; \log |t|+v(\zeta, z)<0\right\}
$$

and if $\Omega^{\prime} \subset \mathbb{C}^{2}$ is the image of $\Omega$ by the projection map $(t, \zeta, z) \longmapsto(t, \zeta)$, then $\Omega^{\prime}=$ $\left\{(t, \zeta) \in \mathbb{C}^{2} ; \log |t|+u(\zeta)<0\right\}$ is not pseudoconvex. However, the minimum property holds true when $v(\zeta, z)$ depends only on $\operatorname{Re} z$ :
(7.5) Theorem ([Kiselman 1978]). Let $\Omega \subset \mathbb{C}^{p} \times \mathbb{C}^{n}$ be a pseudoconvex open set such that each slice

$$
\Omega_{\zeta}=\left\{z \in \mathbb{C}^{n} ;(\zeta, z) \in \Omega\right\}, \quad \zeta \in \mathbb{C}^{p}
$$

is a convex tube $\omega_{\zeta}+\mathbb{R}^{n}$, $\omega_{\zeta} \subset \mathbb{R}^{n}$. For every plurisubharmonic function $v(\zeta, z)$ on $\Omega$ that does not depend on $\operatorname{Im} z$, the function

$$
u(\zeta)=\inf _{z \in \Omega_{\zeta}} v(\zeta, z)
$$

is plurisubharmonic or locally $\equiv-\infty$ on $\Omega^{\prime}=\operatorname{pr}_{\mathbb{C}^{n}}(\Omega)$.
Proof. The hypothesis implies that $v(\zeta, z)$ is convex in $x=\operatorname{Re} z$. In addition, we first assume that $v$ is smooth, plurisubharmonic in $(\zeta, z)$, strictly convex in $x$ and $\lim _{x \rightarrow\{\infty\} \cup \partial \omega_{\zeta}} v(\zeta, x)=+\infty$ for every $\zeta \in \Omega^{\prime}$. Then $x \longmapsto v(\zeta, x)$ has a unique minimum point $x=g(\zeta)$, solution of the equations $\partial v / \partial x_{j}(x, \zeta)=0$. As the matrix $\left(\partial^{2} v / \partial x_{j} \partial x_{k}\right)$ is positive definite, the implicit function theorem shows that $g$ is smooth. Now, if $\mathbb{C} \ni w \longmapsto \zeta_{0}+w a, a \in \mathbb{C}^{n},|w| \leqslant 1$ is a complex disk $\Delta$ contained in $\Omega$, there exists a holomorphic function $f$ on the unit disk, smooth up to the boundary, whose real part solves the Dirichlet problem

$$
\operatorname{Re} f\left(e^{\mathrm{i} \theta}\right)=g\left(\zeta_{0}+e^{\mathrm{i} \theta} a\right)
$$

Since $v\left(\zeta_{0}+w a, f(w)\right)$ is subharmonic in $w$, we get the mean value inequality

$$
v\left(\zeta_{0}, f(0)\right) \leqslant \frac{1}{2 \pi} \int_{0}^{2 \pi} v\left(\zeta_{0}+e^{\mathrm{i} \theta} a, f\left(e^{\mathrm{i} \theta}\right)\right) d \theta=\frac{1}{2 \pi} \int_{\partial \Delta} v(\zeta, g(\zeta)) d \theta
$$

The last equality holds because $\operatorname{Re} f=g$ on $\partial \Delta$ and $v(\zeta, z)=v(\zeta, \operatorname{Re} z)$ by hypothesis. As $u\left(\zeta_{0}\right) \leqslant v\left(\zeta_{0}, f(0)\right)$ and $u(\zeta)=v(\zeta, g(\zeta))$, we see that $u$ satisfies the mean value inequality, thus $u$ is plurisubharmonic.

Now, this result can be extended to arbitrary functions $v$ as follows: let $\psi(\zeta, z) \geqslant 0$ be a continuous plurisubharmonic function on $\Omega$ which is independent of $\operatorname{Im} z$ and is an exhaustion of $\Omega \cap\left(\mathbb{C}^{p} \times \mathbb{R}^{n}\right)$, e.g.

$$
\psi(\zeta, z)=\max \left\{|\zeta|^{2}+|\operatorname{Re} z|^{2},-\log \delta_{\Omega}(\zeta, z)\right\}
$$

There is slowly increasing sequence $C_{j} \rightarrow+\infty$ such that each function $\psi_{j}=\left(C_{j}-\psi \star\right.$ $\left.\rho_{1 / j}\right)^{-1}$ is an "exhaustion" of a pseudoconvex open set $\Omega_{j} \subset \subset \Omega$ whose slices are convex tubes and such that $d\left(\Omega_{j}, C \Omega\right)>2 / j$. Then

$$
v_{j}(\zeta, z)=v \star \rho_{1 / j}(\zeta, z)+\frac{1}{j}|\operatorname{Re} z|^{2}+\psi_{j}(\zeta, z)
$$

is a decreasing sequence of plurisubharmonic functions on $\Omega_{j}$ satisfying our previous conditions. As $v=\lim v_{j}$, we see that $u=\lim u_{j}$ is plurisubharmonic.
(7.6) Corollary. Let $\Omega \subset \mathbb{C}^{p} \times \mathbb{C}^{n}$ be a pseudoconvex open set such that all slices $\Omega_{\zeta}$, $\zeta \in \mathbb{C}^{p}$, are convex tubes in $\mathbb{C}^{n}$. Then the projection $\Omega^{\prime}$ of $\Omega$ on $\mathbb{C}^{p}$ is pseudoconvex.

Proof. Take $v \in \operatorname{Psh}(\Omega)$ equal to the function $\psi$ defined in the proof of Th. 7.5. Then $u$ is a plurisubharmonic exhaustion of $\Omega^{\prime}$.

## $\S$ 7.C. Levi Form of the Boundary

For an arbitrary domain in $\mathbb{C}^{n}$, we first show that pseudoconvexity is a local property of the boundary.
(7.7) Theorem. Let $\Omega \subset \mathbb{C}^{n}$ be an open subset such that every point $z_{0} \in \partial \Omega$ has a neighborhood $V$ such that $\Omega \cap V$ is pseudoconvex. Then $\Omega$ is pseudoconvex.

Proof. As $d(z, \complement \Omega)$ coincides with $d(z, \complement(\Omega \cap V))$ in a neighborhood of $z_{0}$, we see that there exists a neighborhood $U$ of $\partial \Omega$ such that $-\log d(z, \complement \Omega)$ is plurisubharmonic on $\Omega \cap U$. Choose a convex increasing function $\chi$ such that

$$
\chi(r)>\sup _{(\Omega \backslash U) \cap \bar{B}(0, r)}-\log d(z, \mathrm{C} \Omega), \quad \forall r \geqslant 0
$$

Then the function

$$
\psi(z)=\max \{\chi(|z|),-\log d(z,\lceil\Omega)\}
$$

coincides with $\chi(|z|)$ in a neighborhood of $\Omega \backslash U$. Therefore $\psi \in \operatorname{Psh}(\Omega)$, and $\psi$ is clearly an exhaustion.

Now, we give a geometric characterization of the pseudoconvexity property when $\partial \Omega$ is of class $C^{2}$. Let $\rho \in C^{2}(\bar{\Omega})$ be a defining function of $\Omega$, i.e. a function such that

$$
\begin{equation*}
\rho<0 \text { on } \Omega, \quad \rho=0 \text { and } d \rho \neq 0 \text { on } \partial \Omega . \tag{7.9}
\end{equation*}
$$

The holomorphic tangent space to $\partial \Omega$ is by definition the largest complex subspace which is contained in the tangent space $T_{\partial \Omega}$ to the boundary:

$$
\begin{equation*}
{ }^{h} T_{\partial \Omega}=T_{\partial \Omega} \cap J T_{\partial \Omega} . \tag{7.9}
\end{equation*}
$$

It is easy to see that ${ }^{h} T_{\partial \Omega, z}$ is the complex hyperplane of vectors $\xi \in \mathbb{C}^{n}$ such that

$$
d^{\prime} \rho(z) \cdot \xi=\sum_{1 \leqslant j \leqslant n} \frac{\partial \rho}{\partial z_{j}} \xi_{j}=0 .
$$

The Levi form on ${ }^{h} T_{\partial \Omega}$ is defined at every point $z \in \partial \Omega$ by

$$
\begin{equation*}
L_{\partial \Omega, z}(\xi)=\frac{1}{|\nabla \rho(z)|} \sum_{j, k} \frac{\partial^{2} \rho}{\partial z_{j} \partial \bar{z}_{k}} \xi_{j} \bar{\xi}_{k}, \quad \xi \in{ }^{h} T_{\partial \Omega, z} . \tag{7.10}
\end{equation*}
$$

The Levi form does not depend on the particular choice of $\rho$, as can be seen from the following intrinsic computation of $L_{\partial \Omega}$ (we still denote by $L_{\partial \Omega}$ the associated sesquilinear form).
(7.11) Lemma. Let $\xi, \eta$ be $C^{1}$ vector fields on $\partial \Omega$ with values in ${ }^{h} T_{\partial \Omega}$. Then

$$
\langle[\xi, \eta], J \nu\rangle=4 \operatorname{Im} L_{\partial \Omega}(\xi, \eta)
$$

where $\nu$ is the outward normal unit vector to $\partial \Omega$, [, ] the Lie bracket of vector fields and $\langle$,$\rangle the hermitian inner product.$

Proof. Extend first $\xi, \eta$ as vector fields in a neighborhood of $\partial \Omega$ and set

$$
\xi^{\prime}=\sum \xi_{j} \frac{\partial}{\partial z_{j}}=\frac{1}{2}(\xi-\mathrm{i} J \xi), \quad \eta^{\prime \prime}=\sum \bar{\eta}_{k} \frac{\partial}{\partial \bar{z}_{k}}=\frac{1}{2}(\eta+\mathrm{i} J \eta) .
$$

As $\xi, J \xi, \eta, J \eta$ are tangent to $\partial \Omega$, we get on $\partial \Omega$ :

$$
0=\xi^{\prime} \cdot\left(\eta^{\prime \prime} \cdot \rho\right)+\eta^{\prime \prime} \cdot\left(\xi^{\prime} \cdot \rho\right)=\sum_{1 \leqslant j, k \leqslant n} 2 \frac{\partial^{2} \rho}{\partial z_{j} \partial \bar{z}_{k}} \xi_{j} \bar{\eta}_{k}+\xi_{j} \frac{\partial \bar{\eta}_{k}}{\partial z_{j}} \frac{\partial \rho}{\partial \bar{z}_{k}}+\bar{\eta}_{k} \frac{\partial \xi_{j}}{\partial \bar{z}_{k}} \frac{\partial \rho}{\partial z_{j}}
$$

Since $[\xi, \eta]$ is also tangent to $\partial \Omega$, we have $\operatorname{Re}\langle[\xi, \eta], \nu\rangle=0$, hence $\langle J[\xi, \eta], \nu\rangle$ is real and

$$
\langle[\xi, \eta], J \nu\rangle=-\langle J[\xi, \eta], \nu\rangle=-\frac{1}{|\nabla \rho|}(J[\xi, \eta] \cdot \rho)=-\frac{2}{|\nabla \rho|} \operatorname{Re}\left(J\left[\xi^{\prime}, \eta^{\prime \prime}\right] . \rho\right)
$$

because $J\left[\xi^{\prime}, \eta^{\prime}\right]=i\left[\xi^{\prime}, \eta^{\prime}\right]$ and its conjugate $J\left[\xi^{\prime \prime}, \eta^{\prime \prime}\right]$ are tangent to $\partial \Omega$. We find now

$$
\begin{aligned}
J\left[\xi^{\prime}, \eta^{\prime \prime}\right] & =-\mathrm{i} \sum \xi_{j} \frac{\partial \bar{\eta}_{k}}{\partial z_{j}} \frac{\partial}{\partial \bar{z}_{k}}+\bar{\eta}_{k} \frac{\partial \xi_{j}}{\partial \bar{z}_{k}} \frac{\partial}{\partial \bar{z}_{j}}, \\
\operatorname{Re}\left(J\left[\xi^{\prime}, \eta^{\prime \prime}\right] \cdot \rho\right) & =\operatorname{Im} \sum \xi_{j} \frac{\partial \bar{\eta}_{k}}{\partial z_{j}} \frac{\partial \rho}{\partial \bar{z}_{k}}+\bar{\eta}_{k} \frac{\partial \xi_{j}}{\partial \bar{z}_{k}} \frac{\partial \rho}{\partial z_{j}}=-2 \operatorname{Im} \sum \frac{\partial^{2} \rho}{\partial z_{j} \partial \bar{z}_{k}} \xi_{j} \bar{\eta}_{k}, \\
\langle[\xi, \eta], J \nu\rangle & =\frac{4}{|\nabla \rho|} \operatorname{Im} \sum \frac{\partial^{2} \rho}{\partial z_{j} \partial \bar{z}_{k}} \xi_{j} \bar{\eta}_{k}=4 \operatorname{Im} L_{\partial \Omega}(\xi, \eta) .
\end{aligned}
$$

(7.12) Theorem. An open subset $\Omega \subset \mathbb{C}^{n}$ with $C^{2}$ boundary is pseudoconvex if and only if the Levi form $L_{\partial \Omega}$ is semipositive at every point of $\partial \Omega$.

Proof. Set $\delta(z)=d(z, \complement \Omega), z \in \bar{\Omega}$. Then $\rho=-\delta$ is $C^{2}$ near $\partial \Omega$ and satisfies (7.9). If $\Omega$ is pseudoconvex, the plurisubharmonicity of $-\log (-\rho)$ means that for all $z \in \Omega$ near $\partial \Omega$ and all $\xi \in \mathbb{C}^{n}$ one has

$$
\sum_{1 \leqslant j, k \leqslant n}\left(\frac{1}{|\rho|} \frac{\partial^{2} \rho}{\partial z_{j} \partial \bar{z}_{k}}+\frac{1}{\rho^{2}} \frac{\partial \rho}{\partial z_{j}} \frac{\partial \rho}{\partial \bar{z}_{k}}\right) \xi_{j} \bar{\xi}_{k} \geqslant 0 .
$$

Hence $\sum\left(\partial^{2} \rho / \partial z_{j} \partial \bar{z}_{k}\right) \xi_{j} \bar{\xi}_{k} \geqslant 0$ if $\sum\left(\partial \rho / \partial z_{j}\right) \xi_{j}=0$, and an easy argument shows that this is also true at the limit on $\partial \Omega$.

Conversely, if $\Omega$ is not pseudoconvex, Th. 7.2 and 7.7 show that $-\log \delta$ is not plurisubharmonic in any neighborhood of $\partial \Omega$. Hence there exists $\xi \in \mathbb{C}^{n}$ such that

$$
c=\left(\frac{\partial^{2}}{\partial t \partial \bar{t}} \log \delta(z+t \xi)\right)_{\mid t=0}>0
$$

for some $z$ in the neighborhood of $\partial \Omega$ where $\delta \in C^{2}$. By Taylor's formula, we have

$$
\log \delta(z+t \xi)=\log \delta(z)+\operatorname{Re}\left(a t+b t^{2}\right)+c|t|^{2}+\mathrm{o}\left(|t|^{2}\right)
$$

with $a, b \in \mathbb{C}$. Now, choose $z_{0} \in \partial \Omega$ such that $\delta(z)=\left|z-z_{0}\right|$ and set

$$
h(t)=z+t \xi+e^{a t+b t^{2}}\left(z_{0}-z\right), \quad t \in \mathbb{C} .
$$

Then we get $h(0)=z_{0}$ and

$$
\begin{aligned}
\delta(h(t)) & \geqslant \delta(z+t \xi)-\delta(z)\left|e^{a t+b t^{2}}\right| \\
& \geqslant \delta(z)\left|e^{a t+b t^{2}}\right|\left(e^{c|t|^{2} / 2}-1\right) \geqslant \delta(z) c|t|^{2} / 3
\end{aligned}
$$

when $|t|$ is sufficiently small. Since $\delta(h(0))=\delta\left(z_{0}\right)=0$, we obtain at $t=0$ :

$$
\begin{gathered}
\frac{\partial}{\partial t} \delta(h(t))=\sum \frac{\partial \delta}{\partial z_{j}}\left(z_{0}\right) h_{j}^{\prime}(0)=0 \\
\frac{\partial^{2}}{\partial t \partial \bar{t}} \delta(h(t))=\sum \frac{\partial^{2} \delta}{\partial z_{j} \partial \bar{z}_{k}}\left(z_{0}\right) h_{j}^{\prime}(0) \overline{h_{k}^{\prime}(0)}>0
\end{gathered}
$$

hence $h^{\prime}(0) \in{ }^{h} T_{\partial \Omega, z_{0}}$ and $L_{\partial \Omega, z_{0}}\left(h^{\prime}(0)\right)<0$.
(7.13) Definition. The boundary $\partial \Omega$ is said to be weakly (resp. strongly) pseudoconvex if $L_{\partial \Omega}$ is semipositive (resp. positive definite) on $\partial \Omega$. The boundary is said to be Levi flat if $L_{\partial \Omega} \equiv 0$.
(7.14) Remark. Lemma 7.11 shows that $\partial \Omega$ is Levi flat if and only if the subbundle ${ }^{h} T_{\partial \Omega} \subset T_{\partial \Omega}$ is integrable (i.e. stable under the Lie bracket). Assume that $\partial \Omega$ is of class $\mathscr{C}^{k}, k \geqslant 2$. Then ${ }^{h} T_{\partial \Omega}$ is of class $C^{k-1}$. By Frobenius' theorem, the integrability condition implies that ${ }^{h} T_{\partial \Omega}$ is the tangent bundle to a $\mathscr{C}^{k}$ foliation of $\partial \Omega$ whose leaves have real dimension $2 n-2$. But the leaves themselves must be complex analytic since ${ }^{h} T_{\partial \Omega}$ is a complex vector space (cf. Lemma 7.15 below). Therefore $\partial \Omega$ is Levi flat if and only if it is foliated by complex analytic hypersurfaces.
(7.15) Lemma. Let $Y$ be a $C^{1}$-submanifold of a complex analytic manifold $X$. If the tangent space $T_{Y, x}$ is a complex subspace of $T_{X, x}$ at every point $x \in Y$, then $Y$ is complex analytic.

Proof. Let $x_{0} \in Y$. Select holomorphic coordinates $\left(z_{1}, \ldots, z_{n}\right)$ on $X$ centered at $x_{0}$ such that $T_{Y, x_{0}}$ is spanned by $\partial / \partial z_{1}, \ldots, \partial / \partial z_{p}$. Then there exists a neighborhood $U=U^{\prime} \times U^{\prime \prime}$ of $x_{0}$ such that $Y \cap U$ is a graph

$$
z^{\prime \prime}=h\left(z^{\prime}\right), \quad z^{\prime}=\left(z_{1}, \ldots, z_{p}\right) \in U^{\prime}, \quad z^{\prime \prime}=\left(z_{p+1}, \ldots, z_{n}\right)
$$

with $h \in C^{1}\left(U^{\prime}\right)$ and $d h(0)=0$. The differential of $h$ at $z^{\prime}$ is the composite of the projection of $\mathbb{C}^{p} \times\{0\}$ on $T_{Y,\left(z^{\prime}, h\left(z^{\prime}\right)\right)}$ along $\{0\} \times \mathbb{C}^{n-p}$ and of the second projection $\mathbb{C}^{n} \rightarrow \mathbb{C}^{n-p}$. Hence $d h\left(z^{\prime}\right)$ is $\mathbb{C}$-linear at every point and $h$ is holomorphic.

## § 8. Exercises

## $\S$ 8.1. Let $\Omega \subset \mathbb{C}^{n}$ be an open set such that

$$
z \in \Omega, \quad \lambda \in \mathbb{C}, \quad|\lambda| \leqslant 1 \Longrightarrow \lambda z \in \Omega .
$$

Show that $\Omega$ is a union of polydisks of center 0 (with arbitrary linear changes of coordinates) and infer that the space of polynomials $\mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$ is dense in $\mathcal{O}(\Omega)$ for the topology of uniform convergence on compact subsets and in $\mathscr{O}(\Omega) \cap C^{0}(\bar{\Omega})$ for the topology of uniform convergence on $\bar{\Omega}$.
Hint: consider the Taylor expansion of a function $f \in \mathscr{O}(\Omega)$ at the origin, writing it as a series of homogeneous polynomials. To deal with the case of $\bar{\Omega}$, first apply a dilation to $f$.
$\oint$ 8.2. Let $B \subset \mathbb{C}^{n}$ be the unit euclidean ball, $S=\partial B$ and $f \in \mathscr{O}(B) \cap C^{0}(\bar{B})$. Our goal is to check the following Cauchy formula:

$$
f(w)=\frac{1}{\sigma_{2 n-1}} \int_{S} \frac{f(z)}{(1-\langle w, z\rangle)^{n}} d \sigma(z)
$$

a) By means of a unitary transformation and Exercise 8.1, reduce the question to the case when $w=\left(w_{1}, 0, \ldots, 0\right)$ and $f(z)$ is a monomial $z^{\alpha}$.
b) Show that the integral $\int_{B} z^{\alpha} \bar{z}_{1}^{k} d \lambda(z)$ vanishes unless $\alpha=(k, 0, \ldots, 0)$. Compute the value of the remaining integral by the Fubini theorem, as well as the integrals $\int_{S} z^{\alpha} \bar{z}_{1}^{k} d \sigma(z)$.
c) Prove the formula by a suitable power series expansion.
$\S$ 8.3. A current $T \in \mathscr{D}_{p}^{\prime}(M)$ is said to be normal if both $T$ and $d T$ are of order zero, i.e. have measure coefficients.
a) If $T$ is normal and has support contained in a $C^{1}$ submanifold $Y \subset M$, show that there exists a normal current $\Theta$ on $Y$ such that $T=j_{\star} \Theta$, where $j: Y \longrightarrow M$ is the inclusion.
Hint: if $x_{1}=\ldots=x_{q}=0$ are equations of $Y$ in a coordinate system $\left(x_{1}, \ldots, x_{n}\right)$, observe that $x_{j} T=x_{j} d T=0$ for $1 \leqslant j \leqslant q$ and infer that $d x_{1} \wedge \ldots \wedge d x_{q}$ can be factorized in all terms of $T$.
b) What happens if $p>\operatorname{dim} Y$ ?
c) Are a) and b) valid when the normality assumption is dropped ?
$\oint$ 8.4. Let $T=\sum_{1 \leqslant j \leqslant n} T_{j} d \bar{z}_{j}$ be a closed current of bidegree $(0,1)$ with compact support in $\mathbb{C}^{n}$ such that $d^{\prime \prime} T=0$.
a) Show that the partial convolution $S=\left(1 / \pi z_{1}\right) \star_{1} T_{1}$ is a solution of the equation $d^{\prime \prime} S=T$.
b) Let $K=\operatorname{Supp} T$. If $n \geqslant 2$, show that $S$ has support in the compact set $\widetilde{K}$ equal to the union of $K$ and of all bounded components of $\mathbb{C}^{n} \backslash K$.
Hint: observe that $S$ is holomorphic on $\mathbb{C}^{n} \backslash K$ and that $S$ vanishes for $\left|z_{2}\right|+\ldots+\left|z_{n}\right|$ large.
$\oint$ 8.5. Alternative proof of the Dolbeault-Grothendieck lemma. Let $v=\sum_{|J|=q} v_{J} d \bar{z}_{J}, q \geqslant 1$, be a smooth form of bidegree $(0, q)$ on a polydisk $\Omega=D(0, R) \subset \mathbb{C}^{n}$, such that $d^{\prime \prime} v=0$, and let $\omega=D(0, r) \subset \subset \omega$. Let $k$ be the smallest integer such that the monomials $d \bar{z}_{J}$ appearing in $v$ only involve $d \bar{z}_{1}, \ldots, d \bar{z}_{k}$. Prove by induction on $k$ that the equation $d^{\prime \prime} u=v$ can be solved on $\omega$.
Hint: set $v=f \wedge d \bar{z}_{k}+g$ where $f, g$ only involve $d \bar{z}_{1}, \ldots, d \bar{z}_{k-1}$. Then consider $v-d^{\prime \prime} F$ where

$$
F=\sum_{|J|=q-1} F_{J} d \bar{z}_{J}, \quad F_{J}(z)=\left(\psi\left(z_{k}\right) f_{J}\right) \star_{k}\left(\frac{1}{\pi z_{k}}\right),
$$

where $\star_{k}$ denotes the partial convolution with respect to $z_{k}, \psi\left(z_{k}\right)$ is a cut-off function equal to 1 on $D\left(0, r_{k}+\varepsilon\right)$ and $f=\sum_{|J|=q-1} f_{J} d \bar{z}_{J}$.
$\S$ 8.6. Construct locally bounded non continuous subharmonic functions on $\mathbb{C}$.
Hint: consider $e^{u}$ where $u(z)=\sum_{j \geqslant 1} 2^{-j} \log |z-1 / j|$.
$\oint$ 8.7. Let $\omega$ be an open subset of $\mathbb{R}^{n}, n \geqslant 2$, and $u$ a subharmonic function which is not locally $-\infty$.
a) For every open set $\omega \subset \subset \Omega$, show that there is a positive measure $\mu$ with support in $\bar{\omega}$ and a harmonic function $h$ on $\omega$ such that $u=N \star \mu+h$ on $\omega$.
b) Use this representation to prove the following properties: $u \in L_{\text {loc }}^{p}$ for all $p<n /(n-2)$ and $\partial u / \partial x_{j} \in L_{\mathrm{loc}}^{p}$ for all $p<n /(n-1)$.
$\S$ 8.8. Show that a connected open set $\Omega \subset \mathbb{R}^{n}$ is convex if and only if $\Omega$ has a locally convex exhaustion function $\varphi$.
Hint: to show the sufficiency, take a path $\gamma:[0,1] \rightarrow \Omega$ joining two arbitrary points $a, b \in \Omega$ and consider the restriction of $\varphi$ to $\left[a, \gamma\left(t_{0}\right)\right] \cap \Omega$ where $t_{0}$ is the supremum of all $t$ such that $[a, \gamma(u)] \subset \Omega$ for $u \in[0, t]$.
$\S$ 8.9. Let $\left.r_{1}, r_{2} \in\right] 1,+\infty[$. Consider the compact set

$$
K=\left\{\left|z_{1}\right| \leqslant r_{1},\left|z_{2}\right| \leqslant 1\right\} \cup\left\{\left|z_{1}\right| \leqslant 1,\left|z_{2}\right| \leqslant r_{2}\right\} \subset \mathbb{C}^{2}
$$

Show that the holomorphic hull of $K$ in $\mathbb{C}^{2}$ is

$$
\widehat{K}=\left\{\left|z_{1}\right| \leqslant r_{1},\left|z_{2}\right| \leqslant r_{2},\left|z_{1}\right|^{1 / \log r_{1}}\left|z_{2}\right|^{1 / \log r_{2}} \leqslant e\right\} .
$$

Hint: to show that $\widehat{K}$ is contained in this set, consider all holomorphic monomials $f\left(z_{1}, z_{2}\right)=z_{1}^{\alpha_{1}} z_{2}^{\alpha_{2}}$. To show the converse inclusion, apply the maximum principle to the domain $\left|z_{1}\right| \leqslant r_{1},\left|z_{2}\right| \leqslant r_{2}$ on suitably chosen Riemann surfaces $z_{1}^{\alpha_{1}} z_{2}^{\alpha_{2}}=\lambda$.
$\S$ 8.10. Compute the rank of the Levi form of the ellipsoid $\left|z_{1}\right|^{2}+\left|z_{3}\right|^{4}+\left|z_{3}\right|^{6}<1$ at every point of the boundary.
$\S$ 8.11. Let $X$ be a complex manifold and let $u(z)=\sum_{j \in \mathbb{N}}\left|f_{j}\right|^{2}, f_{j} \in \mathscr{O}(X)$, be a series converging uniformly on every compact subset of $X$. Prove that the limit is real analytic and that the series remains uniformly convergent by taking derivatives term by term.
Hint: since the problem is local, take $X=B(0, r)$, a ball in $\mathbb{C}^{n}$. Let $g_{j}(z)=\overline{g_{j}(\bar{z})}$ be the conjugate function of $f_{j}$ and let $U(z, w)=\sum_{j \in \mathbb{N}} f_{j}(z) g_{j}(w)$ on $X \times X$. Using the Cauchy-Schwarz inequality, show that this series of holomorphic functions is uniformly convergent on every compact subset of $X \times X$.
$\S$ 8.12. Let $\Omega \subset \mathbb{C}^{n}$ be a bounded open set with $C^{2}$ boundary.
a) Let $a \in \partial \Omega$ be a given point. Let $e_{n}$ be the outward normal vector to $T_{\partial \Omega, a},\left(e_{1}, \ldots, e_{n-1}\right)$ an orthonormal basis of ${ }^{h} T_{a}(\partial \Omega)$ in which the Levi form is diagonal and $\left(z_{1}, \ldots, z_{n}\right)$ the associated linear coordinates centered at $a$. Show that there is a neighborhood $V$ of $a$ such that $\partial \Omega \cap V$ is the graph $\operatorname{Re} z_{n}=-\varphi\left(z_{1}, \ldots, z_{n-1}, \operatorname{Im} z_{n}\right)$ of a function $\varphi$ such that $\varphi(z)=O\left(|z|^{2}\right)$ and the matrix $\partial^{2} \varphi / \partial z_{j} \partial \bar{z}_{k}(0), 1 \leqslant j, k \leqslant n-1$ is diagonal.
b) Show that there exist local analytic coordinates $w_{1}=z_{1}, \ldots, w_{n-1}=z_{n-1}, w_{n}=z_{n}+\sum c_{j k} z_{j} z_{k}$ on a neighborhood $V^{\prime}$ of $a=0$ such that

$$
\Omega \cap V^{\prime}=V^{\prime} \cap\left\{\operatorname{Re} w_{n}+\sum_{1 \leqslant j \leqslant n} \lambda_{j}\left|w_{j}\right|^{2}+\mathrm{o}\left(|w|^{2}\right)<0\right\}, \quad \lambda_{j} \in \mathbb{R}
$$

and that $\lambda_{n}$ can be assigned to any given value by a suitable choice of the coordinates.
Hint: Consider the Taylor expansion of order 2 of the defining function $\rho(z)=\left(\operatorname{Re} z_{n}+\varphi(z)\right)(1+$ $\operatorname{Re} \sum c_{j} z_{j}$ ) where $c_{j} \in \mathbb{C}$ are chosen in a suitable way.
c) Prove that $\partial \Omega$ is strongly pseudoconvex at $a$ if and only if there is a neighborhood $U$ of $a$ and a biholomorphism $\Phi$ of $U$ onto some open set of $\mathbb{C}^{n}$ such that $\Phi(\Omega \cap U)$ is strongly convex.
d) Assume that the Levi form of $\partial \Omega$ is not semipositive. Show that all holomorphic functions $f \in \mathscr{O}(\Omega)$ extend to some (fixed) neighborhood of $a$.
Hint: assume for example $\lambda_{1}<0$. For $\varepsilon>0$ small, show that $\Omega$ contains the Hartogs figure

$$
\begin{aligned}
\left\{\varepsilon / 2<\left|w_{1}\right|<\varepsilon\right\} & \times\left\{\left|w_{j}\right|<\varepsilon^{2}\right\}_{1<j<n} \times\left\{\left|w_{n}\right|<\varepsilon^{3 / 2}, \operatorname{Re} w_{n}<\varepsilon^{3}\right\} \cup \\
\left\{\left|w_{1}\right|<\varepsilon\right\} & \times\left\{\left|w_{j}\right|<\varepsilon^{2}\right\}_{1<j<n} \times\left\{\left|w_{n}\right|<\varepsilon^{3 / 2}, \operatorname{Re} w_{n}<-\varepsilon^{2}\right\}
\end{aligned}
$$

$\S$ 8.13. Let $\Omega \subset \mathbb{C}^{n}$ be a bounded open set with $C^{2}$ boundary and $\rho \in C^{2}(\Omega, \mathbb{R})$ such that $\rho<0$ on $\Omega, \rho=0$ and $d \rho \neq 0$ on $\partial \Omega$. Let $f \in C^{1}(\partial \Omega, \mathbb{C})$ be a function satisfying the tangential Cauchy-Riemann equations

$$
\xi^{\prime \prime} \cdot f=0, \quad \forall \xi \in{ }^{h} T_{\partial \Omega}, \quad \xi^{\prime \prime}=\frac{1}{2}(\xi+\mathrm{i} J \xi)
$$

a) Let $f_{0}$ be a $C^{1}$ extension of $f$ to $\bar{\Omega}$. Show that $d^{\prime \prime} f_{0} \wedge d^{\prime \prime} \rho=0$ on $\partial \Omega$ and infer that $v=\mathbb{1}_{\Omega} d^{\prime \prime} f_{0}$ is a $d^{\prime \prime}$-closed current on $\mathbb{C}^{n}$.
b) Show that the solution $u$ of $d^{\prime \prime} u=v$ provided by Cor. 3.27 is continuous and that $f$ admits an extension $\widetilde{f} \in \mathscr{O}(\Omega) \cap C^{0}(\bar{\Omega})$ if $\partial \Omega$ is connected.
$\oint$ 8.14. Let $\Omega \subset \mathbb{C}^{n}$ be a bounded pseudoconvex domain with $C^{2}$ boundary and let $\delta(z)=d(z, \mathrm{C} \Omega)$ be the euclidean distance to the boundary.
a) Use the plurisubharmonicity of $-\log \delta$ to prove the following fact: for every $\varepsilon>0$ there is a constant $C_{\varepsilon}>0$ such that

$$
\frac{-H \delta_{z}(\xi)}{\delta(z)}+\varepsilon \frac{\left|d^{\prime} \delta_{z} \cdot \xi\right|^{2}}{|\delta(z)|^{2}}+C_{\varepsilon}|\xi|^{2} \geqslant 0
$$

for $\xi \in \mathbb{C}^{n}$ and $z$ near $\partial \Omega$.
b) Set $\psi(z)=-\log \delta(z)+K|z|^{2}$. Show that for $K$ large and $\alpha$ small the function

$$
\rho(z)=-\exp (-\alpha \psi(z))=-\left(e^{-K|z|^{2}} \delta(z)\right)^{\alpha}
$$

is plurisubharmonic.
c) Prove the existence of a plurisubharmonic exhaustion function $u: \Omega \rightarrow\left[-1,0\left[\right.\right.$ of class $C^{2}$ such that $|u(z)|$ has the same order of magnitude as $\delta(z)^{\alpha}$ when $z$ tends to $\partial \Omega$.
Hint: consult [Diederich-Fornaess 1976].
$\oint$ 8.15. Let $\Omega=\omega+\mathrm{i} \mathbb{R}^{n}$ be a connected tube in $\mathbb{C}^{n}$ of base $\omega$.
a) Assume first that $n=2$. Let $T \subset \mathbb{R}^{2}$ be the triangle $x_{1} \geqslant 0, x_{2} \geqslant 0, x_{1}+x_{2} \leqslant 1$, and assume that the two edges $[0,1] \times\{0\}$ and $\{0\} \times[0,1]$ are contained in $\omega$. Show that every holomorphic function $f \in \mathscr{O}(\Omega)$ extends to a neighborhood of $T+\mathrm{i} \mathbb{R}^{2}$.
Hint: let $\pi: \mathbb{C}^{2} \longrightarrow \mathbb{R}^{2}$ be the projection on the real part and $M_{\varepsilon}$ the intersection of $\pi^{-1}((1+\varepsilon) T)$ with the Riemann surface $z_{1}+z_{2}-\frac{\varepsilon}{2}\left(z_{1}^{2}+z_{2}^{2}\right)=1$ (a non degenerate affine conic). Show that $M_{\varepsilon}$ is compact and that

$$
\begin{aligned}
& \pi\left(\partial M_{\varepsilon}\right) \subset([0,1+\varepsilon] \times\{0\}) \cup(\{0\} \times[0,1+\varepsilon]) \subset \omega \\
& \pi\left([0,1] \cdot M_{\varepsilon}\right) \supset T
\end{aligned}
$$

for $\varepsilon$ small. Use the Cauchy formula along $\partial M_{\varepsilon}$ (in some parametrization of the conic) to obtain an extension of $f$ to $[0,1] \cdot M_{\varepsilon}+\mathbb{R}^{n}$.
b) In general, show that every $f \in \mathscr{O}(\Omega)$ extends to the convex hull $\widehat{\Omega}$.

Hint: given $a, b \in \omega$, consider a polygonal line joining $a$ and $b$ and apply a) inductively to obtain an extension along $[a, b]+\mathrm{i} \mathbb{R}^{n}$.
$\S$ 8.16. For each integer $\nu \geqslant 1$, consider the algebraic variety

$$
X_{\nu}=\left\{(z, w, t) \in \mathbb{C}^{3} ; w t=p_{\nu}(z)\right\}, \quad p_{\nu}(z)=\prod_{1 \leqslant k \leqslant \nu}(z-1 / k)
$$

and the map $j_{\nu}: X_{\nu} \rightarrow X_{\nu+1}$ such that

$$
j_{\nu}(z, w, t)=\left(z, w, t\left(z-\frac{1}{\nu+1}\right)\right) .
$$

a) Show that $X_{\nu}$ is a Stein manifold, and that $j_{\nu}$ is an embedding of $X_{\nu}$ onto an open subset of $X_{\nu+1}$.
b) Define $X=\lim \left(X_{\nu}, j_{\nu}\right)$, and let $\pi_{\nu}: X_{\nu} \rightarrow \mathbb{C}^{2}$ be the projection to the first two coordinates. Since $\pi_{\nu+1} \circ j_{\nu}=\pi_{\nu}$, there exists a holomorphic map $\pi: X \rightarrow \mathbb{C}^{2}, \pi=\lim \pi_{\nu}$. Show that

$$
\mathbb{C}^{2} \backslash \pi(X)=\left\{(z, 0) \in \mathbb{C}^{2} ; z \neq 1 / \nu, \forall \nu \in \mathbb{N}, \nu \geqslant 1\right\}
$$

and especially, that $(0,0) \notin \pi(X)$.
c) Consider the compact set

$$
K=\pi^{-1}\left(\left\{(z, w) \in \mathbb{C}^{2} ;|z| \leqslant 1,|w|=1\right\}\right)
$$

By looking at points of the forms $(1 / \nu, w, 0),|w|=1$, show that $\pi^{-1}(1 / \nu, 1 / \nu) \in \widehat{K}_{0(X)}$. Conclude from this that $X$ is not holomorphically convex (this example is due to [Fornaess 1977]).
$\oint$ 8.17. Let $X$ be a complex manifold, and let $\pi: \widetilde{X} \rightarrow X$ be a holomorphic unramified covering of $X$ ( $X$ and $\widetilde{X}$ are assumed to be connected).
a) Let $g$ be a complete riemannian metric on $X$, and let $\widetilde{d}$ be the geodesic distance on $\widetilde{X}$ associated to $\widetilde{g}=\pi^{\star} g$ (see VIII-2.3 for definitions). Show that $\widetilde{g}$ is complete and that $\delta_{0}(x):=\widetilde{d}\left(x, x_{0}\right)$ is a continuous exhaustion function on $\widetilde{X}$, for any given point $x_{0} \in \widetilde{X}$.
b) Let $\left(U_{\alpha}\right)$ be a locally finite covering of $X$ by open balls contained in coordinate open sets, such that all intersections $U_{\alpha} \cap U_{\beta}$ are diffeomorphic to convex open sets (see Lemma IV-6.9). Let $\theta_{\alpha}$ be a partition of unity subordinate to the covering $\left(U_{\alpha}\right)$, and let $\delta_{\varepsilon_{\alpha}}$ be the convolution of $\delta_{0}$ with a regularizing kernel $\rho_{\varepsilon_{\alpha}}$ on each piece of $\pi^{-1}\left(U_{\alpha}\right)$ which is mapped biholomorphically onto $U_{\alpha}$. Finally, set $\delta=\sum\left(\theta_{\alpha} \circ \pi\right) \delta_{\varepsilon_{\alpha}}$. Show that if $\left(\varepsilon_{\alpha}\right)$ is a collection of sufficiently small positive numbers, then $\delta$ is a smooth exhaustion function on $\widetilde{X}$.
c) Using the fact that $\delta_{0}$ is 1-Lipschitz with respect to $\tilde{d}$, show that derivatives $\partial^{|\nu|} \delta(x) / \partial x^{\nu}$ of a given order with respect to coordinates in $U_{\alpha}$ are uniformly bounded in all components of $\pi^{-1}\left(U_{\alpha}\right)$, at least when $x$ lies in the compact subset $\operatorname{Supp} \theta_{\alpha}$. Conclude from this that there exists a positive hermitian form with continuous coefficients on $X$ such that $H \delta \geqslant-\pi^{\star} \gamma$ on $\widetilde{X}$.
d) If $X$ is strongly pseudoconvex, show that $\tilde{X}$ is also strongly pseudoconvex.

Hint: let $\psi$ be a smooth strictly plurisubharmonic exhaustion function on $X$. Show that there exists a smooth convex increasing function $\chi: \mathbb{R} \rightarrow \mathbb{R}$ such that $\delta+\chi \circ \psi$ is strictly plurisubharmonic.

## Chapter II

## Coherent Sheaves and Analytic Spaces


#### Abstract

The chapter starts with rather general and abstract concepts concerning sheaves and ringed spaces. Introduced in the decade 1950-1960 by Leray, Cartan, Serre and Grothendieck, sheaves and ringed spaces have since been recognized as the adequate tools to handle algebraic varieties and analytic spaces in a unified framework. We then concentrate ourselves on the theory of complex analytic functions. The second section is devoted to a proof of the Weierstrass preparation theorem, which is nothing but a division algorithm for holomorphic functions. It is used to derive algebraic properties of the ring $\mathscr{O}_{n}$ of germs of holomorphic functions in $\mathbb{C}^{n}$. Coherent analytic sheaves are then introduced and the fundamental coherence theorem of Oka is proved. Basic properties of analytic sets are investigated in detail: local parametrization theorem, Hilbert's Nullstellensatz, coherence of the ideal sheaf of an analytic set, analyticity of the singular set. The formalism of complex spaces is then developed and gives a natural setting for the proof of more global properties (decomposition into global irreducible components, maximum principle). After a few definitions concerning cycles, divisors and meromorphic functions, we investigate the important notion of normal space and establish the Oka normalization theorem. Next, the Remmert-Stein extension theorem and the Remmert proper mapping theorem on images of analytic sets are proved by means of semi-continuity results on the rank of morphisms. As an application, we give a proof of Chow's theorem asserting that every analytic subset of $\mathbb{P}^{n}$ is algebraic. Finally, the concept of analytic scheme with nilpotent elements is introduced as a generalization of complex spaces, and we discuss the concepts of bimeromorphic maps, modifications and blowing-up.


## § 1. Presheaves and Sheaves

## $\S$ 1.A. Main Definitions.

Sheaves have become a very important tool in analytic or algebraic geometry as well as in algebraic topology. They are especially useful when one wants to relate global properties of an object to its local properties (the latter being usually easier to establish). We first introduce the axioms of presheaves and sheaves in full generality and give some basic examples.
(1.1) Definition. Let $X$ be a topological space. A presheaf $\mathscr{A}$ on $X$ consists of the following data:
a) a collection of non empty sets $\mathscr{A}(U)$ associated with every open set $U \subset X$,
b) a collection of maps $\rho_{U, V}: \mathscr{A}(V) \longrightarrow \mathscr{A}(U)$ defined whenever $U \subset V$ and satisfying the transitivity property
c) $\rho_{U, V} \circ \rho_{V, W}=\rho_{U, W} \quad$ for $U \subset V \subset W, \quad \rho_{U, U}=\operatorname{Id}_{U} \quad$ for every $U$.

The set $\mathscr{A}(U)$ is called the set of sections of the presheaf $\mathscr{A}$ over $U$.

Most often, the presheaf $\mathscr{A}$ is supposed to carry an additional algebraic structure. For instance:
(1.2) Definition. A presheaf $\mathscr{A}$ is said to be a presheaf of abelian groups (resp. rings, $R$ modules, algebras) if all sets $\mathscr{A}(U)$ are abelian groups (resp. rings, $R$-modules, algebras) and if the maps $\rho_{U, V}$ are morphisms of these algebraic structures. In this case, we always assume that $\mathcal{A}(\emptyset)=\{0\}$.
(1.3) Example. If we assign to each open set $U \subset X$ the set $\mathscr{C}(U)$ of all real valued continuous functions on $U$ and let $\rho_{U, V}$ be the obvious restriction morphism $\mathscr{C}(V) \rightarrow$ $\mathscr{C}(U)$, then $\mathscr{C}$ is a presheaf of rings on $X$. Similarly if $X$ is differentiable (resp. complex analytic) manifold, there are well defined presheaves of rings $\mathscr{C}^{k}$ of functions of class $\mathscr{C}^{k}$ (resp. ©) of holomorphic functions) on $X$. Because of these examples, the maps $\rho_{U, V}$ in Def. 1.1 are often viewed intuitively as "restriction homomorphisms", although the sets $\mathscr{A}(U)$ are not necessarily sets of functions defined over $U$. For the simplicity of notation we often just write $\rho_{U, V}(f)=f_{\lceil U}$ whenever $f \in \mathscr{A}(V), V \supset U$.

For the above presheaves $\mathscr{C}, \mathscr{C}^{k}, \mathfrak{O}$, the properties of functions under consideration are purely local. As a consequence, these presheaves satisfy the following additional gluing axioms, where $\left(U_{\alpha}\right)$ and $U=\bigcup U_{\alpha}$ are arbitrary open subsets of $X$ :

$$
\text { If } F_{\alpha} \in \mathscr{A}\left(U_{\alpha}\right) \text { are such that } \rho_{U_{\alpha} \cap U_{\beta}, U_{\alpha}}\left(F_{\alpha}\right)=\rho_{U_{\alpha} \cap U_{\beta}, U_{\beta}}\left(F_{\beta}\right)
$$ for all $\alpha, \beta$, there exists $F \in \mathscr{A}(U)$ such that $\rho_{U_{\alpha}, U}(F)=F_{\alpha}$;

$$
\text { If } F, G \in \mathscr{A}(U) \text { and } \rho_{U_{\alpha}, U}(F)=\rho_{U_{\alpha}, U}(G) \text { for all } \alpha \text {, then } F=G \text {; }
$$

in other words, local sections over the sets $U_{\alpha}$ can be glued together if they coincide in the intersections and the resulting section on $U$ is uniquely defined. Not all presheaves satisfy (1.4') and (1.4 ${ }^{\prime \prime}$ ):
(1.5) Example. Let $E$ be an arbitrary set with a distinguished element 0 (e.g. an abelian group, a $R$-module, ...). The constant presheaf $E_{X}$ on $X$ is defined to be $E_{X}(U)=E$ for all $\emptyset \neq U \subset X$ and $E_{X}(\emptyset)=\{0\}$, with restriction maps $\rho_{U, V}=\operatorname{Id}_{E}$ if $\emptyset \neq U \subset V$ and $\rho_{U, V}=0$ if $U=\emptyset$. Then axiom (1.4') is not satisfied if $U$ is the union of two disjoint open sets $U_{1}, U_{2}$ and $E$ contains a non zero element.
(1.6) Definition. A presheaf $\mathscr{A}$ is said to be a sheaf if it satisfies the gluing axioms (1.4') and (1.4").

If $\mathscr{A}, \mathscr{B}$ are presheaves of abelian groups (or of some other algebraic structure) on the same space $X$, a presheaf morphism $\varphi: \mathscr{A} \rightarrow \mathscr{B}$ is a collection of morphisms $\varphi_{U}$ : $\mathscr{A}(U) \rightarrow \mathscr{B}(U)$ commuting with the restriction morphisms, i.e. such that for each pair $U \subset V$ there is a commutative diagram

$$
\begin{array}{cr}
\mathscr{A}(V) \xrightarrow{\varphi_{V}} \mathscr{B}(V) \\
\rho_{U, V}^{\mathscr{B}} \downarrow & \downarrow \rho_{U, V}^{\mathscr{A}} \\
\mathscr{A}(U) \xrightarrow{\varphi_{U}} & \mathscr{B}(U) .
\end{array}
$$

We say that $\mathscr{A}$ is a subpresheaf of $\mathscr{B}$ in the case where $\varphi_{U}: \mathscr{A}(U) \subset \mathscr{B}(U)$ is the inclusion morphism; the commutation property then means that $\rho_{U, V}^{\mathscr{B}}(\mathscr{A}(V)) \subset \mathscr{A}(U)$ for all $U$, $V$, and that $\rho_{U, V}^{\mathscr{A}}$ coincides with $\rho_{U, V}^{\mathscr{B}}$ on $\mathscr{A}(V)$. If $\mathscr{A}$ is a subpresheaf of a presheaf $\mathscr{B}$ of abelian groups, there is a presheaf quotient $\mathscr{C}=\mathscr{B} / \mathscr{A}$ defined by $\mathscr{C}(U)=\mathscr{B}(U) / \mathscr{A}(U)$. In a similar way, one defines the presheaf kernel (resp. presheaf image, presheaf cokernel) of a presheaf morphism $\varphi: \mathscr{A} \rightarrow \mathscr{B}$ to be the presheaves

$$
U \mapsto \operatorname{Ker} \varphi_{U}, \quad U \mapsto \operatorname{Im} \varphi_{U}, \quad U \mapsto \operatorname{Coker} \varphi_{U}
$$

The direct sum $\mathscr{A} \oplus \mathscr{B}$ of presheaves of abelian groups $\mathscr{A}, \mathscr{B}$ is the presheaf $U \mapsto \mathscr{A}(U) \oplus$ $\mathscr{B}(U)$, the tensor product $\mathscr{A} \otimes \mathscr{B}$ of presheaves of $R$-modules is $U \mapsto \mathscr{A}(U) \otimes_{R} \mathscr{B}(U)$, etc
(1.7) Remark. The reader should take care of the fact that the presheaf quotient of a sheaf by a subsheaf is not necessarily a sheaf. To give a specific example, let $X=S^{1}$ be the unit circle in $\mathbb{R}^{2}$, let $\mathscr{C}$ be the sheaf of continuous complex valued functions and $\not \mathscr{\not}$ the subsheaf of integral valued continuous functions (i.e. locally constant functions to $\mathbb{Z}$ ). The exponential map

$$
\varphi=\exp (2 \pi \mathrm{i} \bullet): \mathscr{C} \longrightarrow \mathscr{C}^{\star}
$$

is a morphism from $\mathscr{C}$ to the sheaf $\mathscr{C}^{\star}$ of invertible continuous functions, and the kernel of $\varphi$ is precisely $\mathscr{\mathscr { L }}$. However $\varphi_{U}$ is surjective for all $U \neq X$ but maps $\mathscr{C}(X)$ onto the multiplicative subgroup of continuous functions of $\mathscr{C}^{\star}(X)$ of degree 0 . Therefore the quotient presheaf $\mathscr{C} / \mathscr{\mathscr { E }}$ is not isomorphic with $\mathscr{C}^{\star}$, although their groups of sections are the same for all $U \neq X$. Since $\mathscr{C}^{\star}$ is a sheaf, we see that $\mathscr{C} / \mathscr{\not}$ does not satisfy property (1.4').

In order to overcome the difficulty appearing in Example 1.7, it is necessary to introduce a suitable process by which we can produce a sheaf from a presheaf. For this, it is convenient to introduce a slightly modified viewpoint for sheaves.
(1.8) Definition. If $\mathscr{A}$ is a presheaf, we define the set $\widetilde{\mathscr{A}}_{x}$ of germs of $\mathscr{A}$ at a point $x \in X$ to be the abstract inductive limit

$$
\widetilde{\mathscr{A}}_{x}=\underset{U \ni x}{\lim }\left(\mathscr{A}(U), \rho_{U, V}\right) .
$$

More explicitely, $\widetilde{\mathcal{A}}_{x}$ is the set of equivalence classes of elements in the disjoint union $\coprod_{U \ni x} \mathscr{A}(U)$ taken over all open neighborhoods $U$ of $x$, with two elements $F_{1} \in \mathscr{A}\left(U_{1}\right)$, $F_{2} \in \mathscr{A}\left(U_{2}\right)$ being equivalent, $F_{1} \sim F_{2}$, if and only if there is a neighborhood $V \subset U_{1}, U_{2}$ such that $F_{1 \mid V}=F_{2 \mid V}$, i.e., $\rho_{V U_{1}}\left(F_{1}\right)=\rho_{V U_{2}}\left(F_{2}\right)$. The germ of an element $F \in \mathscr{A}(U)$ at a point $x \in U$ will be denoted by $F_{x}$.

Let $\mathscr{A}$ be an arbitrary presheaf. The disjoint union $\widetilde{\mathscr{A}}=\coprod_{x \in X} \widetilde{\mathscr{A}}_{x}$ can be equipped with a natural topology as follows: for every $F \in \mathscr{A}(U)$, we set

$$
\Omega_{F, U}=\left\{F_{x} ; x \in U\right\}
$$

and choose the $\Omega_{F, U}$ to be a basis of the topology of $\tilde{\mathscr{A}}$; note that this family is stable by intersection: $\Omega_{F, U} \cap \Omega_{G, V}=\Omega_{H, W}$ where $W$ is the (open) set of points $x \in U \cap V$ at
which $F_{x}=G_{x}$ and $H=\rho_{W, U}(F)$. The obvious projection map $\pi: \widetilde{\mathscr{A}} \rightarrow X$ which sends $\tilde{\mathscr{A}}_{x}$ to $\{x\}$ is then a local homeomorphism (it is actually a homeomorphism from $\Omega_{F, U}$ onto $U$ ). This leads in a natural way to the following definition:
(1.9) Definition. Let $X$ and $\mathscr{S}$ be topological spaces (not necessarily Hausdorff), and let $\pi: \mathscr{S} \longrightarrow X$ be a mapping such that
a) $\pi$ maps $\mathscr{S}$ onto $X$;
b) $\pi$ is a local homeomorphism, that is, every point in $\mathscr{S}$ has an open neighborhood which is mapped homeomorphically by $\pi$ onto an open subset of $X$.

Then $\mathscr{S}$ is called a sheaf-space on $X$ and $\pi$ is called the projection of $\mathscr{S}$ on $X$. If $x \in X$, then $\mathscr{S}_{x}=\pi^{-1}(x)$ is called the stalk of $\mathscr{I}$ at $x$.

If $Y$ is a subset of $X$, we denote by $\Gamma(Y, \mathscr{F})$ the set of sections of $\mathscr{S}$ on $Y$, i.e. the set of continuous functions $F: Y \rightarrow \mathscr{S}$ such that $\pi \circ F=\operatorname{Id}_{Y}$. It is clear that the presheaf defined by the collection of sets $\mathscr{S}^{\prime}(U):=\Gamma(U, \mathscr{S})$ for all open sets $U \subset X$ together with the restriction maps $\rho_{U, V}$ satisfies axioms (1.4') and (1.4"), hence $\mathscr{S}^{\prime}$ is a sheaf. The set of germs of $\mathscr{S}^{\prime}$ at $x$ is in one-to-one correspondence with the stalk $\mathscr{S}_{x}=\pi^{-1}(x)$, thanks to the local homeomorphism assumption 1.9 b ). This shows that one can associate in a natural way a sheaf $\mathscr{S}^{\prime}$ to every sheaf-space $\mathscr{S}$, and that the sheaf-space $\left(\mathscr{S}^{\prime}\right)^{\sim}$ can be considered to be identical to the original sheaf-space $\mathscr{T}$. Since the assignment $\mathscr{S} \mapsto \mathscr{S}^{\prime}$ from sheaf-spaces to sheaves is an equivalence of categories, we will usually omit the prime sign in the notation of $\mathscr{S}^{\prime}$ and thus use the same symbols for a sheaf-space and its associated sheaf of sections; in a corresponding way, we write $\Gamma(U, \mathscr{Y})=\mathscr{I}(U)$ when $U$ is an open set.

Conversely, given a presheaf $\mathscr{A}$ on $X$, we have an associated sheaf-space $\widetilde{\mathscr{A}}$ and an obvious presheaf morphism

$$
\begin{equation*}
\mathscr{A}(U) \longrightarrow \widetilde{\mathscr{A}}^{\prime}(U)=\Gamma(U, \widetilde{A}), \quad F \longmapsto \widetilde{F}=\left(U \ni x \mapsto F_{x}\right) . \tag{1.10}
\end{equation*}
$$

This morphism is clearly injective if and only if $\mathscr{A}$ satisfies axiom ( $1.4^{\prime \prime}$ ), and it is not difficult to see that $\left(1.4^{\prime}\right)$ and (1.4 $4^{\prime \prime}$ ) together imply surjectivity. Therefore $\mathscr{A} \rightarrow \widetilde{\mathscr{A}}^{\prime}$ is an isomorphism if and only if $\mathscr{A}$ is a sheaf. According to the equivalence of categories between sheaves and sheaf-spaces mentioned above, we will use from now on the same symbol $\widetilde{\mathscr{A}}$ for the sheaf-space and its associated sheaf $\widetilde{\mathscr{A}^{\prime}}$; one says that $\widetilde{\mathscr{A}}$ is the sheaf associated with the presheaf $\mathscr{A}$. If $\mathscr{A}$ itself is a sheaf, we will again identify $\widetilde{\mathscr{A}}$ and $\mathscr{A}$, but we will of course keep the notational difference for a presheaf $\mathscr{A}$ which is not a sheaf.
(1.11) Example. The sheaf associated to the constant presheaf of stalk $E$ over $X$ is the sheaf of locally constant functions $X \rightarrow E$. This sheaf will be denoted merely by $E_{X}$ or $E$ if there is no risk of confusion with the corresponding presheaf. In Example 1.7, we have $\mathscr{Z}=\mathbb{Z}_{X}$ and the sheaf $\left(\mathscr{C} / \mathbb{Z}_{X}\right)^{\sim}$ associated with the quotient presheaf $\mathscr{C} / \mathbb{Z}_{X}$ is isomorphic to $\mathscr{C}^{\star}$ via the exponential map.

In the sequel, we usually work in the category of sheaves rather than in the category of presheaves themselves. For instance, the quotient $\mathscr{B} / \mathscr{A}$ of a sheaf $\mathscr{B}$ by a subsheaf $\mathscr{A}$ generally refers to the sheaf associated with the quotient presheaf: its stalks are equal
to $\mathscr{B}_{x} / \mathscr{A}_{x}$, but a section $G$ of $\mathscr{B} / \mathscr{A}$ over an open set $U$ need not necessarily come from a global section of $\mathscr{B}(U)$; what can be only said is that there is a covering $\left(U_{\alpha}\right)$ of $U$ and local sections $F_{\alpha} \in \mathscr{B}\left(U_{\alpha}\right)$ representing $G_{\mid U_{\alpha}}$ such that $\left(F_{\beta}-F_{\alpha}\right)_{\mid U_{\alpha} \cap U_{\beta}}$ belongs to $\mathscr{A}\left(U_{\alpha} \cap U_{\beta}\right)$. A sheaf morphism $\varphi: \mathscr{A} \rightarrow \mathscr{B}$ is said to be injective (resp. surjective) if the germ morphism $\varphi_{x}: \mathscr{A}_{x} \rightarrow \mathscr{B}_{x}$ is injective (resp. surjective) for every $x \in X$. Let us note again that a surjective sheaf morphism $\varphi$ does not necessarily give rise to surjective morphisms $\varphi_{U}: \mathscr{A}(U) \rightarrow \mathscr{B}(U)$.

## $\S$ 1.B. Direct and Inverse Images of Sheaves

Let $X, Y$ be topological spaces and let $f: X \rightarrow Y$ be a continuous map. If $\mathscr{A}$ is a presheaf on $X$, the direct image $f_{\star} \mathscr{A}$ is the presheaf on $Y$ defined by

$$
\begin{equation*}
f_{\star} \mathscr{A}(U)=\mathscr{A}\left(f^{-1}(U)\right) \tag{1.12}
\end{equation*}
$$

for all open sets $U \subset Y$. When $\mathscr{A}$ is a sheaf, it is clear that $f_{\star} \mathscr{A}$ also satisfies axioms (1.4') and (1.4"), thus $f_{\star} \mathscr{A}$ is a sheaf. Its stalks are given by

$$
\begin{equation*}
\left(f_{\star} \mathscr{A}\right)_{y}=\underset{V \ni y}{\lim _{\longrightarrow}} \mathscr{A}\left(f^{-1}(V)\right) \tag{1.13}
\end{equation*}
$$

where $V$ runs over all open neighborhoods of $y \in Y$.
Now, let $\mathscr{B}$ be a sheaf on $Y$, viewed as a sheaf-space with projection map $\pi: \mathscr{B} \rightarrow Y$. We define the inverse image $f^{-1} \mathscr{B}$ by

$$
\begin{equation*}
f^{-1} \mathscr{B}=\mathscr{B} \times_{Y} X=\{(s, x) \in \mathscr{B} \times X ; \pi(s)=f(x)\} \tag{1.14}
\end{equation*}
$$

with the topology induced by the product topology on $\mathscr{B} \times X$. It is then easy to see that the projection $\pi^{\prime}=\operatorname{pr}_{2}: f^{-1} \mathscr{B} \rightarrow X$ is a local homeomorphism, therefore $f^{-1} \mathscr{B}$ is a sheaf on $X$. By construction, the stalks of $f^{-1} \mathscr{B}$ are

$$
\begin{equation*}
\left(f^{-1} \mathscr{B}\right)_{x}=\mathscr{B}_{f(x)}, \tag{1.15}
\end{equation*}
$$

and the sections $\sigma \in f^{-1} \mathscr{B}(U)$ can be considered as continuous mappings $s: U \rightarrow \mathscr{B}$ such that $\pi \circ \sigma=f$. In particular, any section $s \in \mathscr{B}(V)$ on an open set $V \subset Y$ has a pull-back

$$
\begin{equation*}
f^{\star} s=s \circ f \in f^{-1} \mathscr{B}\left(f^{-1}(V)\right) \tag{1.16}
\end{equation*}
$$

There are always natural sheaf morphisms

$$
\begin{equation*}
f^{-1} f_{\star} \mathscr{A} \longrightarrow \mathscr{A}, \quad \mathscr{B} \longrightarrow f_{\star} f^{-1} \mathscr{B} \tag{1.17}
\end{equation*}
$$

defined as follows. A germ in $\left(f^{-1} f_{\star} \mathscr{A}\right)_{x}=\left(f_{\star} \mathscr{A}\right)_{f(x)}$ is defined by a local section $s \in\left(f_{\star} \mathscr{A}\right)(V)=\mathscr{A}\left(f^{-1}(V)\right)$ for some neighborhood $V$ of $f(x)$; this section can be mapped to the germ $s_{x} \in \mathscr{A}_{x}$. In the opposite direction, the pull-back $f^{\star} s$ of a section $s \in \mathscr{B}(V)$ can be seen by (1.16) as a section of $f_{\star} f^{-1} \mathscr{B}(V)$. It is not difficult to see that these natural morphisms are not isomorphisms in general. For instance, if $f$ is a finite covering map with $q$ sheets and if we take $\mathscr{A}=E_{X}, \mathscr{B}=E_{Y}$ to be constant sheaves, then $f_{\star} E_{X} \simeq E_{Y}^{q}$ and $f^{-1} E_{Y}=E_{X}$, thus $f^{-1} f_{\star} E_{X} \simeq E_{X}^{q}$ and $f_{\star} f^{-1} E_{Y} \simeq E_{Y}^{q}$.

## § 1.C. Ringed Spaces

Many natural geometric structures considered in analytic or algebraic geometry can be described in a convenient way as topological spaces equipped with a suitable "structure sheaf" which, most often, is a sheaf of commutative rings. For instance, a lot of properties of $\mathscr{C}^{k}$ differentiable (resp. real analytic, complex analytic) manifolds can be described in terms of their sheaf of rings $\mathscr{C}_{X}^{k}$ of differentiable functions (resp. $\mathscr{C}_{X}^{\omega}$ of real analytic functions, $\mathscr{O}_{X}$ of holomorphic functions). We first recall a few standard definitions concerning rings, referring to textbooks on algebra for more details (see e.g. [Lang 1965]).
(1.18) Some definitions and conventions about rings. All our rings $R$ are supposed implicitly to have a unit element $1_{R}$ (if $R=\{0\}$, we agree that $1_{R}=0_{R}$ ), and a ring morphism $R \rightarrow R^{\prime}$ is supposed to map $1_{R}$ to $1_{R^{\prime}}$. In the subsequent definitions, we assume that all rings under consideration are commutative.
a) An ideal $I \subset R$ is said to be prime if $x y \in I$ implies $x \in I$ or $y \in I$, i.e., if the quotient ring $R / I$ is entire.
b) An ideal $I \subset R$ is said to be maximal if $I \neq R$ and there are no ideals $J$ such that $I \subsetneq J \subsetneq R$ (equivalently, if the quotient ring $R / I$ is a field).
c) The ring $R$ is said to be a local ring if $R$ has a unique maximal ideal $\mathfrak{m}$ (equivalently, if $R$ has an ideal $\mathfrak{m}$ such that all elements of $R \backslash \mathfrak{m}$ are invertible). Its residual field is defined to be the quotient field $R / \mathfrak{m}$.
d) The ring $R$ is said to be Noetherian if every ideal $I \subset R$ is finitely generated (equivalently, if every increasing sequence of ideals $I_{1} \subset I_{2} \subset \ldots$ is stationary).
e) The radical $\sqrt{I}$ of an ideal $I$ is the set of all elements $x \in R$ such that some power $x^{m}, m \in \mathbb{N}^{\star}$, lies in in $I$. Then $\sqrt{I}$ is again an ideal of $R$.
f) The nilradical $N(R)=\sqrt{\{0\}}$ is the ideal of nilpotent elements of $R$. The ring $R$ is said to be reduced if $N(R)=\{0\}$. Otherwise, its reduction is defined to be the reduced ring $R / N(R)$.

We now introduce the general notion of a ringed space.
(1.19) Definition. A ringed space is a pair $\left(X, \mathscr{R}_{X}\right)$ consisting of a topological space $X$ and of a sheaf of rings $\mathscr{R}_{X}$ on $X$, called the structure sheaf. A morphism

$$
F:\left(X, \mathscr{R}_{X}\right) \rightarrow\left(Y, \mathscr{R}_{Y}\right)
$$

of ringed spaces is a pair $\left(f, F^{\star}\right)$ where $f: X \rightarrow Y$ is a continuous map and

$$
F^{\star}: f^{-1} \mathscr{R}_{Y} \rightarrow \mathscr{R}_{X}, \quad F_{x}^{\star}: \mathscr{R}_{Y, f(x)} \rightarrow \mathscr{R}_{X, x}
$$

a homomorphism of sheaves of rings on $X$, called the comorphism of $F$.
If $F:\left(X, \mathscr{R}_{X}\right) \rightarrow\left(Y, \mathscr{R}_{Y}\right)$ and $G:\left(Y, \mathscr{R}_{Y}\right) \rightarrow\left(Z, \mathscr{R}_{Z}\right)$ are morphisms of ringed spaces, the composite $G \circ F$ is the pair consisting of the map $g \circ f: X \rightarrow Z$ and of the comorphism $(G \circ F)^{\star}=F^{\star} \circ f^{-1} G^{\star}$ :

$$
\begin{align*}
& F^{\star} \circ f^{-1} G^{\star}: f^{-1} g^{-1} \mathscr{R}_{Z} \xrightarrow{f^{-1} G^{\star}} f^{-1} \mathscr{R}_{Y} \xrightarrow{F^{\star}} \mathscr{R}_{X},  \tag{1.20}\\
& F_{x}^{\star} \circ G_{f(x)}^{\star}: \mathscr{R}_{Z, g \circ f(x)} \longrightarrow \mathscr{R}_{Y, f(x)} \longrightarrow \mathscr{R}_{X, x} .
\end{align*}
$$

We say of course that $F$ is an isomorphism of ringed spaces if there exists $G$ such that $G \circ F=\operatorname{Id}_{X}$ and $F \circ G=\operatorname{Id}_{Y}$.

If $\left(X, \mathscr{R}_{X}\right)$ is a ringed space, the nilradical of $\mathscr{R}_{X}$ defines an ideal subsheaf $\mathcal{N}_{X}$ of $\mathscr{R}_{X}$, and the identity map $\operatorname{Id}_{X}: X \rightarrow X$ together with the ring homomorphism $\mathscr{R}_{X} \rightarrow \mathscr{R}_{X} / \mathcal{N}_{X}$ defines a ringed space morphism

$$
\begin{equation*}
\left(X, \mathscr{R}_{X} / \mathcal{N}_{X}\right) \rightarrow\left(X, \mathscr{R}_{X}\right) \tag{1.21}
\end{equation*}
$$

called the reduction morphism. Quite often, the letter $X$ by itself is used to denote the ringed space $\left(X, \mathscr{R}_{X}\right)$; we then denote by $X_{\text {red }}=\left(X, \mathscr{R}_{X} / \mathcal{N}_{X}\right)$ its reduction. The ringed space $X$ is said to be reduced if $\mathcal{N}_{X}=0$, in which case the reduction morphism $X_{\text {red }} \rightarrow X$ is an isomorphism. In all examples considered later on in this book, the structure sheaf $\mathscr{R}_{X}$ will be a sheaf of local rings over some field $k$. The relevant definition is as follows.

## (1.22) Definition.

a) A local ringed space is a ringed space $\left(X, \mathscr{R}_{X}\right)$ such that all stalks $\mathscr{R}_{X, x}$ are local rings. The maximal ideal of $\mathscr{R}_{X, x}$ will be denoted by $\mathfrak{m}_{X, x}$. A morphism $F=\left(f, F^{\star}\right)$ : $\left(X, \mathscr{R}_{X}\right) \rightarrow\left(Y, \mathscr{R}_{Y}\right)$ of local ringed spaces is a morphism of ringed spaces such that $F_{x}^{\star}\left(\mathfrak{m}_{Y, f(x)}\right) \subset \mathfrak{m}_{X, x}$ at any point $x \in X$ (i.e., $F_{x}^{\star}$ is a "local" homomorphism of rings).
b) A local ringed space over a field $k$ is a local ringed space $\left(X, \mathscr{R}_{X}\right)$ such that all rings $\mathscr{R}_{X, x}$ are local $k$-algebras with residual field $\mathscr{R}_{X, x} / \mathfrak{m}_{X, x} \simeq k$. A morphism $F$ between such spaces is supposed to have its comorphism defined by local $k$-homomorphisms $F_{x}^{\star}: \mathscr{R}_{Y, f(x)} \rightarrow \mathscr{R}_{X, x}$.

If $\left(X, \mathscr{R}_{X}\right)$ is a local ringed space over $k$, we can associate to each section $s \in \mathscr{R}_{X}(U)$ a function

$$
\bar{s}: U \rightarrow k, \quad x \mapsto \bar{s}(x) \in k=\mathscr{R}_{X, x} / \mathfrak{m}_{X, x},
$$

and we get a sheaf morphism $\mathscr{R}_{X} \rightarrow \overline{\mathscr{R}}_{X}$ onto a subsheaf of rings $\overline{\mathscr{R}}_{X}$ of the sheaf of functions from $X$ to $k$. We clearly have a factorization

$$
\mathscr{R}_{X} \rightarrow \mathscr{R}_{X} / \mathcal{N}_{X} \rightarrow \overline{\mathscr{R}}_{X},
$$

and thus a corresponding factorization of ringed space morphisms (with $\operatorname{Id}_{X}$ as the underlying set theoretic map)

$$
X_{\text {st-red }} \rightarrow X_{\text {red }} \rightarrow X
$$

where $X_{\text {st-red }}=\left(X, \bar{R}_{X}\right)$ is called the strong reduction of $\left(X, \mathscr{R}_{X}\right)$. It is easy to see that $X_{\text {st-red }}$ is actually a reduced local ringed space over $k$. We say that $X$ is strongly reduced if $\mathscr{R}_{X} \rightarrow \overline{\mathscr{R}}_{X}$ is an isomorphism, that is, if $\mathscr{R}_{X}$ can be identified with a subsheaf of the sheaf of functions $X \rightarrow k$ (in our applications to the theory of algebraic or analytic schemes, the concepts of reduction and strong reduction will actually be the same; in general, these notions differ, see Exercise ??.??). It is important to observe that reduction (resp. strong reduction) is a fonctorial process:
if $F=\left(f, F^{\star}\right):\left(X, \mathscr{R}_{X}\right) \rightarrow\left(Y, \mathscr{R}_{Y}\right)$ is a morphism of ringed spaces (resp. of local ringed spaces over $k$ ), there are natural reductions

$$
\begin{aligned}
F_{\text {red }} & =\left(f, F_{\text {red }}^{\star}\right): X_{\text {red }} \rightarrow Y_{\text {red }}, \quad F_{\text {red }}^{\star}: \mathscr{R}_{Y, f(x)} / \mathcal{N}_{Y, f(x)} \rightarrow \mathscr{R}_{X, x} / \mathcal{N}_{X, x}, \\
F_{\text {st-red }} & =\left(f, f^{\star}\right): X_{\text {st-red }} \rightarrow Y_{\text {st-red }}, \quad f^{\star}: \overline{\mathscr{R}}_{Y, f(x)} \rightarrow \overline{\mathscr{R}}_{X, x}, \quad \bar{s} \mapsto \bar{s} \circ f
\end{aligned}
$$

where $f^{\star}$ is the usual pull-back comorphism associated with $f$. Therefore, if ( $X, \mathscr{R}_{X}$ ) and $\left(Y, \mathscr{R}_{Y}\right)$ are strongly reduced, the morphism $F$ is completely determined by the underlying set-theoretic map $f$. Our first basic examples of (strongly reduced) ringed spaces are the various types of manifolds already defined in Chapter I. The language of ringed spaces provides an equivalent but more elegant and more intrinsic definition.
(1.23) Definition. Let $X$ be a Hausdorff separable topological space. One can define the category of $\mathscr{C}^{k}, k \in \mathbb{N} \cup\{\infty, \omega\}$, differentiable manifolds (resp. complex analytic manifolds) to be the category of reduced local ringed spaces $\left(X, \mathscr{R}_{X}\right)$ over $\mathbb{R}$ (resp. over $\left.\mathbb{C}\right)$, such that every point $x \in X$ has a neighborhood $U$ on which the restriction $\left(U, \mathscr{R}_{X \mid U}\right)$ is isomorphic to a ringed space $\left(\Omega, \mathscr{G}_{\Omega}^{k}\right)$ where $\Omega \subset \mathbb{R}^{n}$ is an open set and $\mathscr{G}_{\Omega}^{k}$ is the sheaf of $\mathscr{G}^{k}$ differentiable functions (resp. $\left(\Omega, \mathscr{O}_{\Omega}\right)$, where $\Omega \subset \mathbb{C}^{n}$ is an open subset, and $\mathscr{O}_{\Omega}$ is the sheaf of holomorphic functions on $\Omega$ ).

We say that the ringed spaces $\left(\Omega, \mathscr{C}_{\Omega}^{k}\right)$ and $\left(\Omega, \mathscr{O}_{\Omega}\right)$ are the models of the category of differentiable (resp. complex analytic) manifolds, and that a general object ( $X, \mathscr{R}_{X}$ ) in the category is locally isomorphic to one of the given model spaces. It is easy to see that the corresponding ringed spaces morphisms are nothing but the usual concepts of differentiable and holomorphic maps.

## § 1.D. Algebraic Varieties over a Field

As a second illustration of the notion of ringed space, we present here a brief introduction to the formalism of algebraic varieties, referring to [Hartshorne 1977] or [EGA 1967] for a much more detailed exposition. Our hope is that the reader who already has some background of analytic or algebraic geometry will find some hints of the strong interconnections between both theories. Beginners are invited to skip this section and proceed directly to the theory of complex analytic sheaves in $\S, 2$. All rings or algebras occurring in this section are supposed to be commutative rings with unit.
§ 1.D.1. Affine Algebraic Sets. Let $k$ be an algebraically closed field of any characteristic. An affine algebraic set is a subset $X \subset k^{N}$ of the affine space $k^{N}$ defined by an arbitrary collection $S \subset k\left[T_{1}, \ldots, T_{N}\right]$ of polynomials, that is,

$$
X=V(S)=\left\{\left(z_{1}, \ldots, z_{N}\right) \in k^{N} ; P\left(z_{1}, \ldots, z_{N}\right)=0, \forall P \in S\right\}
$$

Of course, if $J \subset k\left[T_{1}, \ldots, T_{N}\right]$ is the ideal generated by $S$, then $V(S)=V(J)$. As $k\left[T_{1}, \ldots, T_{N}\right]$ is Noetherian, $J$ is generated by finitely many elements $\left(P_{1}, \ldots, P_{m}\right)$, thus $X=V\left(\left\{P_{1}, \ldots, P_{m}\right\}\right)$ is always defined by finitely many equations. Conversely, for any subset $Y \subset k^{N}$, we consider the ideal $I(Y)$ of $k\left[T_{1}, \ldots, T_{N}\right]$, defined by

$$
I(Y)=\left\{P \in k\left[T_{1}, \ldots, T_{N}\right] ; P(z)=0, \forall z \in Y\right\} .
$$

Of course, if $Y \subset k^{N}$ is an algebraic set, we have $V(I(Y))=Y$. In the opposite direction, we have the following fundamental result.
(1.24) Hilbert's Nullstellensatz (see [Lang 1965]). If $J \subset k\left[T_{1}, \ldots, T_{N}\right]$ is an ideal, then $I(V(J))=\sqrt{J}$.

If $X=V(J) \subset k^{N}$ is an affine algebraic set, we define the (reduced) ring $\mathbb{O}(X)$ of algebraic functions on $X$ to be the set of all functions $X \rightarrow k$ which are restrictions of polynomials, i.e.,

$$
\begin{equation*}
\mathfrak{O}(X)=k\left[T_{1}, \ldots, T_{N}\right] / I(X)=k\left[T_{1}, \ldots, T_{N}\right] / \sqrt{J} . \tag{1.25}
\end{equation*}
$$

This is clearly a reduced $k$-algebra. An (algebraic) morphism of affine algebraic sets $X=V(J) \subset k^{N}, Y=V\left(J^{\prime}\right) \subset k^{N^{\prime}}$ is a map $f: Y \rightarrow X$ which is the restriction of a polynomial map $k^{N^{\prime}}$ tok $k^{N}$. We then get a $k$-algebra homomomorphism

$$
f^{\star}: \mathscr{O}(X) \rightarrow \mathscr{O}(Y), \quad s \mapsto s \circ f
$$

called the comorphism of $f$. In this way, we have defined a contravariant fonctor

$$
\begin{equation*}
X \mapsto \mathscr{O}(X), \quad f \mapsto f^{\star} \tag{1.26}
\end{equation*}
$$

from the category of affine algebraic sets to the category of finitely generated reduced $k$-algebras.

We are going to show the existence of a natural fonctor going in the opposite direction. In fact, let us start with an arbitrary finitely generated algebra $A$ (not necessarily reduced at this moment). For any choice of generators $\left(g_{1}, \ldots, g_{N}\right)$ of $A$ we get a surjective morphism of the polynomial ring $k\left[T_{1}, \ldots, T_{N}\right]$ onto $A$,

$$
k\left[T_{1}, \ldots, T_{N}\right] \rightarrow A, \quad T_{j} \mapsto g_{j}
$$

and thus $A \simeq k\left[T_{1}, \ldots, T_{N}\right] / J$ with the ideal $J$ being the kernel of this morphism. It is well-known that every maximal ideal $\mathfrak{m}$ of $A$ has codimension 1 in $A$ (see [Lang 1965]), so that $\mathfrak{m}$ gives rise to a $k$-algebra homomorphism $A \rightarrow A / \mathfrak{m}=k$. We thus get a bijection

$$
\operatorname{Hom}_{\text {alg }}(A, k) \rightarrow \operatorname{Spm}(A), \quad u \mapsto \operatorname{Ker} u
$$

between the set of $k$-algebra homomorphisms and the set $\operatorname{Spm}(A)$ of maximal ideals of $A$. In fact, if $A=k\left[T_{1}, \ldots, T_{N}\right] / J$, an element $\varphi \in \operatorname{Hom}_{\text {alg }}(A, k)$ is completely determined by the values $z_{j}=\varphi\left(T_{j} \bmod J\right)$, and the corresponding algebra homomorphism $k\left[T_{1}, \ldots, T_{N}\right] \rightarrow k, P \mapsto P\left(z_{1}, \ldots, z_{N}\right)$ can be factorized $\bmod J$ if and only if $z=\left(z_{1}, \ldots, z_{N}\right) \in k^{N}$ satisfies the equations

$$
P\left(z_{1}, \ldots, z_{N}\right)=0, \quad \forall P \in J
$$

We infer from this that

$$
\operatorname{Spm}(A) \simeq V(J)=\left\{\left(z_{1}, \ldots, z_{N}\right) \in k^{N} ; P\left(z_{1}, \ldots, z_{N}\right)=0, \forall P \in J\right\}
$$

can be identified with the affine algebraic set $V(J) \subset k^{N}$. If we are given an algebra homomorphism $\Phi: A \rightarrow B$ of finitely generated $k$-algebras we get a corresponding map $\operatorname{Spm}(\Phi): \operatorname{Spm}(B) \rightarrow \operatorname{Spm}(A)$ described either as

$$
\begin{aligned}
& \operatorname{Spm}(B) \rightarrow \operatorname{Spm}(A), \quad \mathfrak{m} \mapsto \Phi^{-1}(\mathfrak{m}) \quad \text { or } \\
& \operatorname{Hom}_{\mathrm{alg}}(B, k) \rightarrow \operatorname{Hom}_{\mathrm{alg}}(A, k), \quad v \mapsto v \circ \Phi .
\end{aligned}
$$

If $B=k\left[T_{1}^{\prime}, \ldots, T_{N^{\prime}}^{\prime}\right] / J^{\prime}$ and $\operatorname{Spm}(B)=V\left(J^{\prime}\right) \subset k^{N^{\prime}}$, it is easy to see that $\operatorname{Spm}(\Phi)$ : $\operatorname{Spm}(B) \rightarrow \operatorname{Spm}(A)$ is the restriction of the polynomial map

$$
f: k^{N^{\prime}} \rightarrow k^{N}, \quad w \mapsto f(w)=\left(P_{1}(w), \ldots, P_{N}(w)\right)
$$

where $P_{j} \in k\left[T_{1}^{\prime}, \ldots, T_{N^{\prime}}^{\prime}\right]$ are polynomials such that $P_{j}=\Phi\left(T_{j}\right) \bmod J^{\prime}$ in $B$. We have in this way defined a contravariant fonctor

$$
\begin{equation*}
A \mapsto \operatorname{Spm}(A), \quad \Phi \mapsto \operatorname{Spm}(\Phi) \tag{1.27}
\end{equation*}
$$

from the category of finitely generated $k$-algebras to the category of affine algebraic sets.
Since $A=k\left[T_{1}, \ldots, T_{N}\right] / J$ and its reduction $A / N(A)=k\left[T_{1}, \ldots, T_{N}\right] / \sqrt{J}$ give rise to the same algebraic set

$$
V(J)=\operatorname{Spm}(A)=\operatorname{Spm}(A / N(A))=V(\sqrt{J})
$$

we see that the category of affine algebraic sets is actually equivalent to the subcategory of reduced finitely generated $k$-algebras.
(1.28) Example. The simplest example of an affine algebraic set is the affine space

$$
k^{N}=\operatorname{Spm}\left(k\left[T_{1}, \ldots, T_{N}\right]\right),
$$

in particular $\operatorname{Spm}(k)=k^{0}$ is just one point. We agree that $\operatorname{Spm}(\{0\})=\emptyset$ (observe that $V(J)=\emptyset$ when $J$ is the unit ideal in $\left.k\left[T_{1}, \ldots, T_{N}\right]\right)$.
§ 1.D.2. Zariski Topology and Affine Algebraic Schemes. Let $A$ be a finitely generated algebra and $X=\operatorname{Spm}(A)$. To each ideal $\mathfrak{a} \subset A$ we associate the zero variety $V(\mathfrak{a}) \subset X$ which consists of all elements $\mathfrak{m} \in X=\operatorname{Spm}(A)$ such that $\mathfrak{m} \supset \mathfrak{a}$; if

$$
A \simeq k\left[T_{1}, \ldots, T_{N}\right] / J \quad \text { and } \quad X \simeq V(J) \subset k^{N},
$$

then $V(\mathfrak{a})$ can be identified with the zero variety $V\left(J_{\mathfrak{a}}\right) \subset X$ of the inverse image $J_{\mathfrak{a}}$ of $\mathfrak{a}$ in $k\left[T_{1}, \ldots, T_{N}\right]$. For any family $\left(\mathfrak{a}_{\alpha}\right)$ of ideals in $A$ we have

$$
V\left(\sum \mathfrak{a}_{\alpha}\right)=\bigcap V\left(\mathfrak{a}_{\alpha}\right), \quad V\left(\mathfrak{a}_{1}\right) \cup V\left(\mathfrak{a}_{2}\right)=V\left(\mathfrak{a}_{1} \mathfrak{a}_{2}\right),
$$

hence there exists a unique topology on $X$ such that the closed sets consist precisely of all algebraic subsets $(V(\mathfrak{a}))_{\mathfrak{a} \subset A}$ of $X$. This topology is called the Zariski topology. The Zariski topology is almost never Hausdorff (for example, if $X=k$ is the affine line, the open sets are $\emptyset$ and complements of finite sets, thus any two nonempty open sets have nonempty intersection). However, $X$ is a Noetherian space, that is, a topological space in which every decreasing sequence of closed sets is stationary; an equivalent definition is to require that every open set is quasi-compact (from any open covering of an open set, one can extract a finite covering).

We now come to the concept of affine open subsets. For $s \in A$, the open set $D(s)=$ $X \backslash V(s)$ can be given the structure of an affine algebraic variety. In fact, if $A=$ $k\left[T_{1}, \ldots, T_{N}\right] / J$ and $s$ is represented by a polynomial in $k\left[T_{1}, \ldots, T_{N}\right]$, the localized ring
$A[1 / s]$ can be written as $A[1 / s]=k\left[T_{1}, \ldots, T_{N}, T_{N+1}\right] / J_{s}$ where $J_{s}=J\left[T_{N+1}\right]+\left(s T_{N+1}-\right.$ 1), thus

$$
V\left(J_{s}\right)=\{(z, w) \in V(J) \times k ; s(z) w=1\} \simeq V(I) \backslash s^{-1}(0)
$$

and $D(s)$ can be identified with $\operatorname{Spm}(A[1 / s])$. We have $D\left(s_{1}\right) \cap D\left(s_{2}\right)=D\left(s_{1} s_{2}\right)$, and the sets $(D(s))_{s \in A}$ are easily seen to be a basis of the Zariski topology on $X$. The open sets $D(s)$ are called affine open sets. Since the open sets $D(s)$ containing a given point $x \in X$ form a basis of neighborhoods, one can define a sheaf space $\mathbb{0}_{X}$ such that the ring of germs $\mathbb{O}_{X, x}$ is the inductive limit

$$
\mathbb{O}_{X, x}=\underset{D(s) \ni x}{\lim _{\vec{\prime}}} A[1 / s]=\{\text { fractions } p / q ; p, q \in A, q(x) \neq 0\}
$$

This is a local ring with maximal ideal

$$
\mathfrak{m}_{X, x}=\{p / q ; p, q \in A, p(x)=0, q(x) \neq 0\}
$$

and residual field $\mathfrak{O}_{X, x} / \mathfrak{m}_{X, x}=k$. In this way, we get a ringed space $\left(X, \mathfrak{O}_{X}\right)$ over $k$. It is easy to see that $\Gamma\left(X, \mathscr{O}_{X}\right)$ coincides with the finitely generated $k$-algebra $A$. In fact, from the definition of $\mathbb{O}_{X}$, a global section is obtained by gluing together local sections $p_{j} / s_{j}$ on affine open sets $D\left(s_{j}\right)$ with $\bigcup D\left(s_{j}\right)=X, 1 \leqslant j \leqslant m$. This means that the ideal $\mathfrak{a}=\left(s_{1}, \ldots, s_{m}\right) \subset A$ has an empty zero variety $V(\mathfrak{a})$, thus $\mathfrak{a}=A$ and there are elements $u_{j} \in A$ with $\sum u_{j} s_{j}=1$. The compatibility condition $p_{j} / s_{j}=p_{k} / s_{k}$ implies that these elements are induced by

$$
\sum u_{j} p_{j} / \sum u_{j} s_{j}=\sum u_{j} p_{j} \in A
$$

as desired. More generally, since the open sets $D(s)$ are affine, we get

$$
\Gamma\left(D(s), \mathscr{O}_{X}\right)=A[1 / s]
$$

It is easy to see that the ringed space $\left(X, \mathscr{O}_{X}\right)$ is reduced if and only if $A$ itself is reduced; in this case, $X$ is even strongly reduced as Hilbert's Nullstellensatz shows. Otherwise, the reduction $X_{\text {red }}$ can obtained from the reduced algebra $A_{\text {red }}=A / N(A)$.

Ringed spaces $\left(X, \mathscr{O}_{X}\right)$ as above are called affine algebraic schemes over $k$ (although substantially different from the usual definition, our definition can be shown to be equivalent in this special situation; compare with [Hartshorne 1977]); see also Exercise ??.??). The category of affine algebraic schemes is equivalent to the category of finitely generated $k$-algebras (with the arrows reversed).
$\S$ 1.D.3. Algebraic Schemes. Algebraic schemes over $k$ are defined to be ringed spaces over $k$ which are locally isomorphic to affine algebraic schemes, modulo an ad hoc separation condition.
(1.29) Definition. An algebraic scheme over $k$ is a local ringed space $\left(X, \mathcal{O}_{X}\right)$ over $k$ such that
a) $X$ has a finite covering by open sets $U_{\alpha}$ such that $\left(U_{\alpha}, \mathscr{O}_{X \mid U_{\alpha}}\right)$ is isomorphic as a ringed space to an affine algebraic scheme $\left(\operatorname{Spm}\left(A_{\alpha}\right), \mathscr{O}_{\mathrm{Spm}\left(A_{\alpha}\right)}\right)$.
b) $X$ satisfies the algebraic separation axiom, namely the diagonal $\Delta_{X}$ of $X \times X$ is closed for the Zariski topology.

A morphism of algebraic schemes is just a morphism of the underlying local ringed spaces. An (abstract) algebraic variety is the same as a reduced algebraic scheme.

In the above definition, some words of explanation are needed for b), since the product $X \times Y$ of algebraic schemes over $k$ is not the ringed space theoretic product, i.e., the product topological space equipped with the structure sheaf $\operatorname{pr}_{1}^{\star} \mathscr{O}_{X} \otimes_{k} \operatorname{pr}_{2}^{\star} \mathbb{Q}_{Y}$. Instead, we define the product of two affine algebraic schemes $X=\operatorname{Spm}(A)$ and $Y=\operatorname{Spm}(B)$ to be $X \times Y=\operatorname{Spm}\left(A \otimes_{k} B\right)$, equipped with the Zariski topology and the structural sheaf associated with $A \otimes_{k} B$. Notice that the Zariski topology on $X \times Y$ is not the product topology of the Zariski topologies on $X, Y$, as the example $k^{2}=k \times k$ shows; also, the rational function $1 /\left(1-z_{1}-z_{2}\right) \in \mathscr{O}_{k^{2},(0,0)}$ is not in $\mathbb{O}_{k, 0} \otimes_{k} \mathbb{O}_{k, 0}$. In general, if $X, Y$ are written as $X=\bigcup U_{\alpha}$ and $Y=\bigcup V_{\beta}$ with affine open sets $U_{\alpha}, V_{\beta}$, we define $X \times Y$ to be the union of all open affine charts $U_{\alpha} \times V_{\beta}$ with their associated structure sheaves of affine algebraic varieties, the open sets of $X \times Y$ being all unions of open sets in the various charts $U_{\alpha} \times V_{\beta}$. The separation axiom b ) is introduced for the sake of excluding pathological examples such as an affine line $k \amalg\left\{0^{\prime}\right\}$ with the origin changed into a double point.
$\S$ 1.D.4. Subschemes. If $\left(X, \mathscr{O}_{X}\right)$ is an affine algebraic scheme and $A=\Gamma\left(X, \mathscr{0}_{X}\right)$ is the associated algebra, we say that $\left(Y, \mathscr{O}_{Y}\right)$ is a subscheme of $\left(X, \mathscr{O}_{X}\right)$ if there is an ideal $\mathfrak{a}$ of $A$ such that $Y \hookrightarrow X$ is the morphism defined by the algebra morphism $A \rightarrow A / \mathfrak{a}$ as its comorphism. As $\operatorname{Spm}(A / \mathfrak{a}) \rightarrow \operatorname{Spm}(A)$ has for image the set $V(\mathfrak{a})$ of maximal ideals $\mathfrak{m}$ of $A$ containing $\mathfrak{a}$, we see that $Y=V(\mathfrak{a})$ as a set; let us introduce the ideal subsheaf $\mathscr{G}=\mathfrak{a} \mathscr{O}_{X} \subset \mathfrak{O}_{X}$. Since the structural sheaf $\mathscr{O}_{Y}$ is obtained by taken localizations $A / \mathfrak{a}[1 / s]$, it is easy to see that $\mathfrak{O}_{Y}$ coincides with the quotient sheaf $\mathbb{O}_{X} / \mathscr{I}$ restricted to $Y$. Since $\mathfrak{a}$ has finitely many generators, the ideal sheaf $\mathscr{f}$ is locally finitely generated (see § 2 below). This leads to the following definition.
(1.30) Definition. If $\left(X, \mathscr{O}_{X}\right)$ is an algebraic scheme, a (closed) subscheme is an algebraic scheme $\left(Y, \mathscr{O}_{Y}\right)$ such that $Y$ is a Zariski closed subset of $X$, and there is a locally finitely generated ideal subsheaf $\mathscr{G} \subset \mathscr{O}_{X}$ such that $Y=V(\mathscr{G})$ and $\mathscr{O}_{Y}=\left(\mathscr{O}_{X} / \mathscr{F}\right)_{\uparrow Y}$.

If $\left(Y, \mathscr{O}_{Y}\right),\left(Z, \mathscr{O}_{Z}\right)$ are subschemes of $\left(X, \mathscr{O}_{X}\right)$ defined by ideal subsheaves $\mathscr{g}, \mathscr{g}^{\prime} \subset$ $\mathbb{0}_{X}$, there are corresponding subschemes $Y \cap Z$ and $Y \cup Z$ defined as ringed spaces

$$
\left(Y \cap Z, \mathscr{O}_{X} /\left(\mathscr{G}+\mathscr{g}^{\prime}\right)\right), \quad\left(Y \cup Z, \mathscr{O}_{X} / \mathscr{G} \mathscr{g}^{\prime}\right)
$$

$\S$ 1.D.5. Projective Algebraic Varieties. A very important subcategory of the category of algebraic varieties is provided by projective algebraic varieties. Let $\mathbb{P}_{k}^{N}$ be the projective $N$-space, that is, the set $k^{N+1} \backslash\{0\} / k^{\star}$ of equivalence classes of $(N+1)$ tuples $\left(z_{0}, \ldots, z_{N}\right) \in k^{N+1} \backslash\{0\}$ under the equivalence relation given by $\left(z_{0}, \ldots, z_{N}\right) \sim$ $\lambda\left(z_{0}, \ldots, z_{N}\right), \lambda \in k^{\star}$. The corresponding element of $\mathbb{P}_{k}^{N}$ will be denoted $\left[z_{0}: z_{1}: \ldots\right.$ : $\left.z_{N}\right]$. It is clear that $\mathbb{P}_{N}^{k}$ can be covered by the $(N+1)$ affine charts $U_{\alpha}, 0 \leqslant \alpha \leqslant N$, such that

$$
U_{\alpha}=\left\{\left[z_{0}: z_{1}: \ldots: z_{N}\right] \in \mathbb{P}_{k}^{N} \quad z_{\alpha} \neq 0\right\} .
$$

The set $U_{\alpha}$ can be identified with the affine $N$-space $k^{N}$ by the map

$$
U_{\alpha} \rightarrow k^{N}, \quad\left[z_{0}: z_{1}: \ldots: z_{N}\right] \mapsto\left(\frac{z_{0}}{z_{\alpha}}, \frac{z_{1}}{z_{\alpha}}, \ldots, \frac{z_{\alpha-1}}{z_{\alpha}}, \frac{z_{\alpha+1}}{z_{\alpha}}, \ldots, \frac{z_{N}}{z_{\alpha}}\right)
$$

With this identification, $\mathscr{O}\left(U_{\alpha}\right)$ is the algebra of homogeneous rational functions of degree 0 in $z_{0}, \ldots, z_{N}$ which have just a power of $z_{\alpha}$ in their denominator. It is easy to see that the structure sheaves $\mathfrak{O}_{U_{\alpha}}$ and $\mathscr{O}_{U_{\beta}}$ coincide in the intersections $U_{\alpha} \cap U_{\beta}$; they can be glued together to define an algebraic variety structure $\left(\mathbb{P}_{k}^{N}, \mathscr{O}_{\mathbb{P}^{N}}\right)$, such that $\mathbb{O}_{\mathbb{P}^{N},[z]}$ consists of all homogeneous rational functions $p / q$ of degree 0 (i.e., $\operatorname{deg} p=\operatorname{deg} q$ ), such that $q(z) \neq 0$.
(1.30) Definition. An algebraic scheme or variety $\left(X, \mathscr{O}_{X}\right)$ is said to be projective if it is isomorphic to a closed subscheme of some projective space $\left(\mathbb{P}_{k}^{N}, \mathscr{O}_{\mathbb{P}^{N}}\right)$.

We now indicate a standard way of constructing projective schemes. Let $S$ be a collection of homogeneous polynomials $P_{j} \in k\left[z_{0}, \ldots, z_{N}\right]$, of degree $d_{j} \in \mathbb{N}$. We define an associated projective algebraic set

$$
\widetilde{V}(S)=\left\{\left[z_{0}: \ldots: z_{N}\right] \in \mathbb{P}_{k}^{N} ; P(z)=0, \forall P \in S\right\}
$$

Let $J$ be the homogeneous ideal of $k\left[z_{0}, \ldots, z_{N}\right]$ generated by $S$ (recall that an ideal $J$ is said to be homogeneous if $J=\bigoplus J_{m}$ is the direct sum of its homogeneous components, or equivalently, if $J$ is generated by homogeneous elements). We have an associated graded algebra

$$
B=k\left[z_{0}, \ldots, z_{N}\right] / J=\bigoplus B_{m}, \quad B_{m}=k\left[z_{0}, \ldots, z_{N}\right]_{m} / J_{m}
$$

such that $B$ is generated by $B_{1}$ and $B_{m}$ is a finite dimensional vector space over $k$ for each $k$. This is enough to construct the desired scheme structure on $\widetilde{V}(J):=\bigcap \widetilde{V}\left(J_{m}\right)$, as we see in the next subsection.
§ 1.D.6. Projective Scheme Associated with a Graded Algebra. Let us start with a reduced graded $k$-algebra

$$
B=\bigoplus_{m \in \mathbb{N}} B_{m}
$$

such that $B$ is generated by $B_{0}$ and $B_{1}$ as an algebra, and $B_{0}, B_{1}$ are finite dimensional vector spaces over $k$ (it then follows that $B$ is finitely generated and that all $B_{m}$ are finite dimensional vector spaces). Given $s \in B_{m}, m>0$, we define a $k$-algebra $A_{s}$ to be the ring of all fractions of homogeneous degree 0 with a power of $s$ as their denominator, i.e.,

$$
\begin{equation*}
A_{s}=\left\{p / s^{d} ; p \in B_{d m}, d \in \mathbb{N}\right\} \tag{1.31}
\end{equation*}
$$

Since $A_{s}$ is generated by $\frac{1}{s} B_{1}^{m}$ over $B_{0}, A_{s}$ is a finitely generated algebra. We define $U_{s}=\operatorname{Spm}\left(A_{s}\right)$ to be the associated affine algebraic variety. For $s \in B_{m}$ and $s^{\prime} \in B_{m^{\prime}}$, we clearly have algebra homomorphisms

$$
A_{s} \rightarrow A_{s s^{\prime}}, \quad A_{s^{\prime}} \rightarrow A_{s s^{\prime}}
$$

since $A_{s s^{\prime}}$ is the algebra of all 0-homogeneous fractions with powers of $s$ and $s^{\prime}$ in the denominator. As $A_{s s^{\prime}}$ is the same as the localized ring $A_{s}\left[s^{m^{\prime}} / s^{\prime m}\right]$, we see that $U_{s s^{\prime}}$ can be identified with an affine open set in $U_{s}$, and we thus get canonical injections

$$
U_{s s^{\prime}} \hookrightarrow U_{s}, \quad U_{s s^{\prime}} \hookrightarrow U_{s^{\prime}} .
$$

(1.32) Definition. If $B=\bigoplus_{m \in \mathbb{N}} B_{m}$ is a reduced graded algebra generated by its finite dimensional vector subspaces $B_{0}$ and $B_{1}$, we associate an algebraic scheme $\left(X, \mathscr{O}_{X}\right)=$ $\operatorname{Proj}(B)$ as follows. To each finitely generated algebra $A_{s}=\left\{p / s^{d} ; p \in B_{d m}, d \in \mathbb{N}\right\}$ we associate an affine algebraic variety $U_{s}=\operatorname{Spm}\left(A_{s}\right)$. We let $X$ be the union of all open charts $U_{s}$ with the identifications $U_{s} \cap U_{s^{\prime}}=U_{s s^{\prime}}$; then the collection $\left(U_{s}\right)$ is a basis of the topology of $X$, and $\mathfrak{O}_{X}$ is the unique sheaf of local $k$-algebras such that $\Gamma\left(U_{s}, \mathbb{Q}_{X}\right)=A_{s}$ for each $U_{s}$.

The following proposition shows that only finitely many open charts are actually needed to describe $X$ (as required in Def. 1.29 a)).
(1.33) Lemma. If $s_{0}, \ldots, s_{N}$ is a basis of $B_{1}$, then $\operatorname{Proj}(B)=\bigcup_{0 \leqslant j \leqslant N} U_{s_{j}}$.

Proof. In fact, if $x \in X$ is contained in a chart $U_{s}$ for some $s \in B_{m}$, then $U_{s}=$ $\operatorname{Spm}\left(A_{s}\right) \neq \emptyset$, and therefore $A_{s} \neq\{0\}$. As $A_{s}$ is generated by $\frac{1}{s} B_{1}^{m}$, we can find a fraction $f=s_{j_{1}} \ldots s_{j_{m}} / s$ representing an element $f \in \mathscr{O}\left(U_{s}\right)$ such that $f(x) \neq 0$. Then $x \in U_{s} \backslash f^{-1}(0)$, and $U_{s} \backslash f^{-1}(0)=\operatorname{Spm}\left(A_{s}[1 / f]\right)=U_{s} \cap U_{s_{j_{1}}} \cap \ldots \cap U_{s_{j_{m}}}$. In particular $x \in U_{s_{j_{1}}}$.
(1.34) Example. One can consider the projective space $\mathbb{P}_{k}^{N}$ to be the algebraic scheme

$$
\mathbb{P}_{k}^{N}=\operatorname{Proj}\left(k\left[T_{0}, \ldots, T_{N}\right]\right) .
$$

The Proj construction is fonctorial in the following sense: if we have a graded homomorphism $\Phi: B \rightarrow B^{\prime}$ (i.e. an algebra homomorphism such that $\Phi\left(B_{m}\right) \subset B_{m}^{\prime}$, then there are corresponding morphisms $A_{s} \rightarrow A_{\Phi(s)}^{\prime}, U_{\Phi(s)}^{\prime} \rightarrow U_{s}$, and we thus find a scheme morphism

$$
F: \operatorname{Proj}\left(B^{\prime}\right) \rightarrow \operatorname{Proj}(B) .
$$

Also, since $p / s^{d}=p s^{l} / s^{d+l}$, the algebras $A_{s}$ depend only on components $B_{m}$ of large degree, and we have $A_{s}=A_{s^{l}}$. It follows easily that there is a canonical isomorphism

$$
\operatorname{Proj}(B) \simeq \operatorname{Proj}\left(\bigoplus_{m} B_{l m}\right) .
$$

Similarly, we may if we wish change a finite number of components $B_{m}$ without affecting $\operatorname{Proj}(B)$. In particular, we may alway assume that $B_{0}=k 1_{B}$. By selecting finitely many generators $g_{0}, \ldots, g_{N}$ in $B_{1}$, we then find a surjective graded homomorphism $k\left[T_{0}, \ldots, T_{N}\right] \rightarrow B$, thus $B \simeq k\left[T_{0}, \ldots, T_{N}\right] / J$ for some graded ideal $J \subset B$. The algebra homomorphism $k\left[T_{0}, \ldots, T_{N}\right] \rightarrow B$ therefore yields a scheme embedding $\operatorname{Proj}(B) \rightarrow \mathbb{P}^{N}$ onto $V(J)$.

We will not pursue further the study of algebraic varieties from this point of view ; in fact we are mostly interested in the case $k=\mathbb{C}$, and algebraic varieties over $\mathbb{C}$ are a special case of the more general concept of complex analytic space.

## §2. The Local Ring of Germs of Analytic Functions

## §2.A. The Weierstrass Preparation Theorem

Our first goal is to establish a basic factorization and division theorem for analytic functions of several variables, which is essentially due to Weierstrass. We follow here a simple proof given by C.L. Siegel, based on a clever use of the Cauchy formula. Let $g$ be a holomorphic function defined on a neighborhood of 0 in $\mathbb{C}^{n}, g \not \equiv 0$. There exists a dense set of vectors $v \in \mathbb{C}^{n} \backslash\{0\}$ such that the function $\mathbb{C} \ni t \longmapsto g(t v)$ is not identically zero. In fact the Taylor series of $g$ at the origin can be written

$$
g(t v)=\sum_{k=0}^{+\infty} \frac{1}{k!} t^{k} g^{(k)}(v)
$$

where $g^{(k)}$ is a homogeneous polynomial of degree $k$ on $\mathbb{C}^{n}$ and $g^{\left(k_{0}\right)} \not \equiv 0$ for some index $k_{0}$. Thus it suffices to select $v$ such that $g^{\left(k_{0}\right)}(v) \neq 0$. After a change of coordinates, we may assume that $v=(0, \ldots, 0,1)$. Let $s$ be the vanishing order of $z_{n} \longmapsto g\left(0, \ldots, 0, z_{n}\right)$ at $z_{n}=0$. There exists $r_{n}>0$ such that $g\left(0, \ldots, 0, z_{n}\right) \neq 0$ when $0<\left|z_{n}\right| \leqslant r_{n}$. By continuity of $g$ and compactness of the circle $\left|z_{n}\right|=r_{n}$, there exists $r^{\prime}>0$ and $\varepsilon>0$ such that

$$
g\left(z^{\prime}, z_{n}\right) \neq 0 \quad \text { for } z^{\prime} \in \mathbb{C}^{n-1}, \quad\left|z^{\prime}\right| \leqslant r^{\prime}, \quad r_{n}-\varepsilon \leqslant\left|z_{n}\right| \leqslant r_{n}+\varepsilon
$$

For every integer $k \in \mathbb{N}$, let us consider the integral

$$
S_{k}\left(z^{\prime}\right)=\frac{1}{2 \pi \mathrm{i}} \int_{\left|z_{n}\right|=r_{n}} \frac{1}{g\left(z^{\prime}, z_{n}\right)} \frac{\partial g}{\partial z_{n}}\left(z^{\prime}, z_{n}\right) z_{n}^{k} d z_{n}
$$

Then $S_{k}$ is holomorphic in a neighborhood of $\left|z^{\prime}\right| \leqslant r^{\prime}$. Rouché's theorem shows that $S_{0}\left(z^{\prime}\right)$ is the number of roots $z_{n}$ of $g\left(z^{\prime}, z_{n}\right)=0$ in the disk $\left|z_{n}\right|<r_{n}$, thus by continuity $S_{0}\left(z^{\prime}\right)$ must be a constant $s$. Let us denote by $w_{1}\left(z^{\prime}\right), \ldots, w_{s}\left(z^{\prime}\right)$ these roots, counted with multiplicity. By definition of $r_{n}$, we have $w_{1}(0)=\ldots=w_{s}(0)=0$, and by the choice of $r^{\prime}, \varepsilon$ we have $\left|w_{j}\left(z^{\prime}\right)\right|<r_{n}-\varepsilon$ for $\left|z^{\prime}\right| \leqslant r^{\prime}$. The Cauchy residue formula yields

$$
S_{k}\left(z^{\prime}\right)=\sum_{j=1}^{s} w_{j}\left(z^{\prime}\right)^{k}
$$

Newton's formula shows that the elementary symmetric function $c_{k}\left(z^{\prime}\right)$ of degree $k$ in $w_{1}\left(z^{\prime}\right), \ldots, w_{s}\left(z^{\prime}\right)$ is a polynomial in $S_{1}\left(z^{\prime}\right), \ldots, S_{k}\left(z^{\prime}\right)$. Hence $c_{k}\left(z^{\prime}\right)$ is holomorphic in a neighborhood of $\left|z^{\prime}\right| \leqslant r^{\prime}$. Let us set

$$
P\left(z^{\prime}, z_{n}\right)=z_{n}^{s}-c_{1}\left(z^{\prime}\right) z_{n}^{s-1}+\cdots+(-1)^{s} c_{s}\left(z^{\prime}\right)=\prod_{j=1}^{s}\left(z_{n}-w_{j}\left(z^{\prime}\right)\right)
$$

For $\left|z^{\prime}\right| \leqslant r^{\prime}$, the quotient $f=g / P$ (resp. $f=P / g$ ) is holomorphic in $z_{n}$ on the disk $\left|z_{n}\right|<r_{n}+\varepsilon$, because $g$ and $P$ have the same zeros with the same multiplicities, and $f\left(z^{\prime}, z_{n}\right)$ is holomorphic in $z^{\prime}$ for $r_{n}-\varepsilon \leqslant\left|z_{n}\right| \leqslant r_{n}+\varepsilon$. Therefore

$$
f\left(z^{\prime}, z_{n}\right)=\frac{1}{2 \pi \mathrm{i}} \int_{\left|w_{n}\right|=r_{n}+\varepsilon} \frac{f\left(z^{\prime}, w_{n}\right) d w_{n}}{w_{n}-z_{n}}
$$

is holomorphic in $z$ on a neighborhood of the closed polydisk $\bar{\Delta}\left(r^{\prime}, r_{n}\right)=\left\{\left|z^{\prime}\right| \leqslant r^{\prime}\right\} \times$ $\left\{\left|z_{n}\right| \leqslant r_{n}\right\}$. Thus $g / P$ is invertible and we obtain:
(2.1) Weierstrass preparation theorem. Let $g$ be holomorphic on a neighborhood of 0 in $\mathbb{C}^{n}$, such that $g\left(0, z_{n}\right) / z_{n}^{s}$ has a not zero finite limit at $z_{n}=0$. With the above choice of $r^{\prime}$ and $r_{n}$, one can write $g(z)=u(z) P\left(z^{\prime}, z_{n}\right)$ where $u$ is an invertible holomorphic function in a neighborhood of the polydisk $\bar{\Delta}\left(r^{\prime}, r_{n}\right)$, and $P$ is a Weierstrass polynomial in $z_{n}$, that is, a polynomial of the form

$$
P\left(z^{\prime}, z_{n}\right)=z_{n}^{s}+a_{1}\left(z^{\prime}\right) z_{n}^{s-1}+\cdots+a_{s}\left(z^{\prime}\right), \quad a_{k}(0)=0
$$

with holomorphic coefficients $a_{k}\left(z^{\prime}\right)$ on a neighborhood of $\left|z^{\prime}\right| \leqslant r^{\prime}$ in $\mathbb{C}^{n-1}$.
(2.2) Remark. If $g$ vanishes at order $m$ at 0 and $v \in \mathbb{C}^{n} \backslash\{0\}$ is selected such that $g^{(m)}(v) \neq 0$, then $s=m$ and $P$ must also vanish at order $m$ at 0 . In that case, the coefficients $a_{k}\left(z^{\prime}\right)$ are such that $a_{k}\left(z^{\prime}\right)=O\left(\left|z^{\prime}\right|^{k}\right), 1 \leqslant k \leqslant s$.
(2.3) Weierstrass division theorem. Every bounded holomorphic function $f$ on $\Delta=$ $\Delta\left(r^{\prime}, r_{n}\right)$ can be represented in the form

$$
\begin{equation*}
f(z)=g(z) q(z)+R\left(z^{\prime}, z_{n}\right), \tag{2.4}
\end{equation*}
$$

where $q$ and $R$ are analytic in $\Delta, R\left(z^{\prime}, z_{n}\right)$ is a polynomial of degree $\leqslant s-1$ in $z_{n}$, and

$$
\begin{equation*}
\sup _{\Delta}|q| \leqslant C \sup _{\Delta}|f|, \quad \sup _{\Delta}|R| \leqslant C \sup _{\Delta}|f| \tag{2.5}
\end{equation*}
$$

for some constant $C \geqslant 0$ independent of $f$. The representation (2.4) is unique.
Proof (Siegel). It is sufficient to prove the result when $g(z)=P\left(z^{\prime}, z_{n}\right)$ is a Weierstrass polynomial.

Let us first prove the uniqueness. If $f=P q_{1}+R_{1}=P q_{2}+R_{2}$, then

$$
P\left(q_{2}-q_{1}\right)+\left(R_{2}-R_{1}\right)=0 .
$$

It follows that the $s$ roots $z_{n}$ of $P\left(z^{\prime}, \bullet\right)=0$ are zeros of $R_{2}-R_{1}$. Since $\operatorname{deg}_{z_{n}}\left(R_{2}-R_{1}\right) \leqslant$ $s-1$, we must have $R_{2}-R_{1} \equiv 0$, thus $q_{2}-q_{1} \equiv 0$.

In order to prove the existence of $(q, R)$, we set

$$
q\left(z^{\prime}, z_{n}\right)=\lim _{\varepsilon \rightarrow 0+} \frac{1}{2 \pi \mathrm{i}} \int_{\left|w_{n}\right|=r_{n}-\varepsilon} \frac{f\left(z^{\prime}, w_{n}\right)}{P\left(z^{\prime}, w_{n}\right)\left(w_{n}-z_{n}\right)} d w_{n}, \quad z \in \Delta ;
$$

observe that the integral does not depend on $\varepsilon$ when $\varepsilon<r_{n}-\left|z_{n}\right|$ is small enough. Then $q$ is holomorphic on $\Delta$. The function $R=f-P q$ is also holomorphic on $\Delta$ and

$$
R(z)=\lim _{\varepsilon \rightarrow 0+} \frac{1}{2 \pi \mathrm{i}} \int_{\left|w_{n}\right|=r_{n}-\varepsilon} \frac{f\left(z^{\prime}, w_{n}\right)}{P\left(z^{\prime}, w_{n}\right)}\left[\frac{P\left(z^{\prime}, w_{n}\right)-P\left(z^{\prime}, z_{n}\right)}{\left(w_{n}-z_{n}\right)}\right] d w_{n}
$$

The expression in brackets has the form

$$
\left[\left(w_{n}^{s}-z_{n}^{s}\right)+\sum_{j=1}^{s} a_{j}\left(z^{\prime}\right)\left(w_{n}^{s-j}-z_{n}^{s-j}\right)\right] /\left(w_{n}-z_{n}\right)
$$

hence is a polynomial in $z_{n}$ of degree $\leqslant s-1$ with coefficients that are holomorphic functions of $z^{\prime}$. Thus we have the asserted decomposition $f=P q+R$ and

$$
\sup _{\Delta}|R| \leqslant C_{1} \sup _{\Delta}|f|
$$

where $C_{1}$ depends on bounds for the $a_{j}\left(z^{\prime}\right)$ and on $\mu=\min \left|P\left(z^{\prime}, z_{n}\right)\right|$ on the compact set $\left\{\left|z^{\prime}\right| \leqslant r^{\prime}\right\} \times\left\{\left|z_{n}\right|=r_{n}\right\}$. By the maximum principle applied to $q=(f-R) / P$ on each disk $\left\{z^{\prime}\right\} \times\left\{\left|z_{n}\right|<r_{n}-\varepsilon\right\}$, we easily get

$$
\sup _{\Delta}|q| \leqslant \mu^{-1}\left(1+C_{1}\right) \sup _{\Delta}|f| .
$$

## § 2.B. Algebraic Properties of the Ring $\mathbb{O}_{n}$

We give here important applications of the Weierstrass preparation theorem to the study of the ring of germs of holomorphic functions in $\mathbb{C}^{n}$.
(2.6) Notation. We let $\mathfrak{O}_{n}$ be the ring of germs of holomorphic functions on $\mathbb{C}^{n}$ at 0 . Alternatively, $\mathscr{O}_{n}$ can be identified with the ring $\mathbb{C}\left\{z_{1}, \ldots, z_{n}\right\}$ of convergent power series in $z_{1}, \ldots, z_{n}$.
(2.7) Theorem. The ring $\mathfrak{O}_{n}$ is Noetherian, i.e. every ideal $\mathscr{I}$ of $\mathscr{O}_{n}$ is finitely generated.

Proof. By induction on $n$. For $n=1, \mathscr{O}_{n}$ is principal: every ideal $\mathscr{I} \neq\{0\}$ is generated by $z^{s}$, where $s$ is the minimum of the vanishing orders at 0 of the non zero elements of $\mathscr{I}$. Let $n \geqslant 2$ and $\mathscr{I} \subset \mathscr{O}_{n}, \mathscr{I} \neq\{0\}$. After a change of variables, we may assume that $\mathscr{I}$ contains a Weierstrass polynomial $P\left(z^{\prime}, z_{n}\right)$. For every $f \in \mathscr{I}$, the Weierstrass division theorem yields

$$
f(z)=P\left(z^{\prime}, z_{n}\right) q(z)+R\left(z^{\prime}, z_{n}\right), \quad R\left(z^{\prime}, z_{n}\right)=\sum_{k=0}^{s-1} c_{k}\left(z^{\prime}\right) z_{n}^{k}
$$

and we have $R \in \mathscr{I}$. Let us consider the set $\mathscr{M}$ of coefficients $\left(c_{0}, \ldots, c_{s-1}\right)$ in $\mathscr{O}_{n-1}^{\oplus s}$ corresponding to the polynomials $R\left(z^{\prime}, z_{n}\right)$ which belong to $\mathscr{I}$. Then $\mathcal{M}$ is a $\mathscr{O}_{n-1^{-}}$ submodule of $\mathscr{O}_{n-1}^{\oplus s}$. By the induction hypothesis $\mathscr{O}_{n-1}$ is Noetherian; furthermore, every submodule of a finitely generated module over a Noetherian ring is finitely generated ([Lang 1965], Chapter VI). Therefore $\mathscr{M}$ is finitely generated, and $\mathscr{I}$ is generated by $P$ and by polynomials $R_{1}, \ldots, R_{N}$ associated with a finite set of generators of $\mathscr{M}$.

Before going further, we need two lemmas which relate the algebraic properties of $\mathbb{O}_{n}$ to those of the polynomial ring $\mathscr{O}_{n-1}\left[z_{n}\right]$.
(2.8) Lemma. Let $P, F \in \mathscr{O}_{n-1}\left[z_{n}\right]$ where $P$ is a Weierstrass polynomial. If $P$ divides $F$ in $\mathfrak{O}_{n}$, then $P$ divides $F$ in $\mathfrak{O}_{n-1}\left[z_{n}\right]$.

Proof. Assume that $F\left(z^{\prime}, z_{n}\right)=P\left(z^{\prime}, z_{n}\right) h(z), h \in \mathscr{O}_{n}$. The standard division algorithm of $F$ by $P$ in $\mathbb{O}_{n-1}\left[z_{n}\right]$ yields

$$
F=P Q+R, \quad Q, R \in \mathscr{O}_{n-1}\left[z_{n}\right], \quad \operatorname{deg} R<\operatorname{deg} P .
$$

The uniqueness part of Th. 2.3 implies $h(z)=Q\left(z^{\prime}, z_{n}\right)$ and $R \equiv 0$.
(2.9) Lemma. Let $P\left(z^{\prime}, z_{n}\right)$ be a Weierstrass polynomial.
a) If $P=P_{1} \ldots P_{N}$ with $P_{j} \in \mathscr{O}_{n-1}\left[z_{n}\right]$, then, up to invertible elements of $\mathscr{O}_{n-1}$, all $P_{j}$ are Weierstrass polynomials.
b) $P\left(z^{\prime}, z_{n}\right)$ is irreducible in $\mathscr{O}_{n}$ if and only if it is irreducible in $\mathscr{O}_{n-1}\left[z_{n}\right]$.

Proof. a) Assume that $P=P_{1} \ldots P_{N}$ with polynomials $P_{j} \in \mathbb{O}_{n-1}\left[z_{n}\right]$ of respective degrees $s_{j}, \sum_{1 \leqslant j \leqslant N} s_{j}=s$. The product of the leading coefficients of $P_{1}, \ldots, P_{N}$ in $\mathscr{O}_{n-1}$ is equal to 1 ; after normalizing these polynomials, we may assume that $P_{1}, \ldots, P_{N}$ are unitary and $s_{j}>0$ for all $j$. Then

$$
P\left(0, z_{n}\right)=z_{n}^{s}=P_{1}\left(0, z_{n}\right) \ldots P_{N}\left(0, z_{n}\right)
$$

hence $P_{j}\left(0, z_{n}\right)=z_{n}^{s_{j}}$ and therefore $P_{j}$ is a Weierstrass polynomial.
b) Set $s=\operatorname{deg} P$ and $P\left(0, z_{n}\right)=z_{n}^{s}$. Assume that $P$ is reducible in $\mathscr{C}_{n}$, with $P\left(z^{\prime}, z_{n}\right)=$ $g_{1}(z) g_{2}(z)$ for non invertible elements $g_{1}, g_{2} \in \mathscr{O}_{n}$. Then $g_{1}\left(0, z_{n}\right)$ and $g_{2}\left(0, z_{n}\right)$ have vanishing orders $s_{1}, s_{2}>0$ with $s_{1}+s_{2}=s$, and

$$
g_{j}=u_{j} P_{j}, \quad \operatorname{deg} P_{j}=s_{j}, \quad j=1,2,
$$

where $P_{j}$ is a Weierstrass polynomial and $u_{j} \in \mathscr{O}_{n}$ is invertible. Therefore $P_{1} P_{2}=u P$ for an invertible germ $u \in \mathbb{O}_{n}$. Lemma 2.8 shows that $P$ divides $P_{1} P_{2}$ in $\mathscr{O}_{n-1}\left[z_{n}\right]$; since $P_{1}, P_{2}$ are unitary and $s=s_{1}+s_{2}$, we get $P=P_{1} P_{2}$, hence $P$ is reducible in $\mathfrak{O}_{n-1}\left[z_{n}\right]$. The converse implication is obvious from a).
(2.10) Theorem. $\mathscr{O}_{n}$ is a factorial ring, i.e. $\mathfrak{O}_{n}$ is entire and:
a) every non zero germ $f \in \mathbb{O}_{n}$ admits a factorization $f=f_{1} \ldots f_{N}$ in irreducible elements;
b) the factorization is unique up to invertible elements.

Proof. The existence part a) follows from Lemma 2.9 if we take $f$ to be a Weierstrass polynomial and $f=f_{1} \ldots f_{N}$ be a decomposition of maximal length $N$ into polynomials of positive degree. In order to prove the uniqueness, it is sufficient to verify the following statement:
bTf $g$ is an irreducible element that divides a product $f_{1} f_{2}$, then $g$ divides either $f_{1}$ or $f_{2}$.

By Th. 2.1, we may assume that $f_{1}, f_{2}, g$ are Weierstrass polynomials in $z_{n}$. Then $g$ is irreducible and divides $f_{1} f_{2}$ in $\mathscr{O}_{n-1}\left[z_{n}\right]$ thanks to Lemmas 2.8 and 2.9 b ). By induction on $n$, we may assume that $\mathscr{O}_{n-1}$ is factorial. The standard Gauss lemma ([Lang 1965], Chapter V) says that the polynomial ring $A[T]$ is factorial if the ring $A$ is factorial. Hence $\mathfrak{O}_{n-1}\left[z_{n}\right]$ is factorial by induction and thus $g$ must divide $f_{1}$ or $f_{2}$ in $\mathscr{O}_{n-1}\left[z_{n}\right]$.
(2.11) Theorem. If $f, g \in \mathscr{O}_{n}$ are relatively prime, then the germs $f_{z}, g_{z}$ at every point $z \in \mathbb{C}^{n}$ near 0 are again relatively prime.

Proof. One may assume that $f=P, g=Q$ are Weierstrass polynomials. Let us recall that unitary polynomials $P, Q \in \mathscr{A}[X](\mathscr{A}=$ a factorial ring $)$ are relatively prime if and only if their resultant $R \in \mathscr{A}$ is non zero. Then the resultant $R\left(z^{\prime}\right) \in \mathscr{O}_{n-1}$ of $P\left(z^{\prime}, z_{n}\right)$ and $Q\left(z^{\prime}, z_{n}\right)$ has a non zero germ at 0 . Therefore the germ $R_{z^{\prime}}$ at points $z^{\prime} \in \mathbb{C}^{n-1}$ near 0 is also non zero.

## § 3. Coherent Sheaves

## § 3.1. Locally Free Sheaves and Vector Bundles

Section 9 will greatly develope this philosophy. Before introducing the more general notion of a coherent sheaf, we discuss the notion of locally free sheaves over a sheaf a ring. All rings occurring in the sequel are supposed to be commutative with unit (the non commutative case is also of considerable interest, e.g. in view of the theory of $\mathscr{D}$-modules, but this subject is beyond the scope of the present book).
(3.1) Definition. Let $\mathscr{A}$ be a sheaf of rings on a topological space $X$ and let $\mathscr{S}$ a sheaf of modules over $\mathscr{A}$ (or briefly a $\mathscr{A}$-module). Then $\mathscr{S}$ is said to be locally free of rank $r$ over $\mathscr{A}$, if $\mathscr{S}$ is locally isomorphic to $\mathscr{A}^{\oplus r}$ on a neighborhood of every point, i.e. for every $x_{0} \in X$ one can find a neighborhood $\Omega$ and sections $F_{1}, \ldots, F_{r} \in \mathscr{S}(\Omega)$ such that the sheaf homomorphism

$$
F: \mathscr{A}_{\upharpoonright \Omega}^{\oplus r} \longrightarrow \mathscr{S}_{\upharpoonright \Omega}, \quad \mathscr{A}_{x}^{\oplus r} \ni\left(w_{1}, \ldots, w_{r}\right) \longmapsto \sum_{1 \leqslant j \leqslant r} w_{j} F_{j, x} \in \mathscr{S}_{x}
$$

is an isomorphism.

By definition, if $\mathscr{S}$ is locally free, there is a covering $\left(U_{\alpha}\right)_{\alpha \in I}$ by open sets on which $\mathscr{S}$ admits free generators $F_{\alpha}^{1}, \ldots, F_{\alpha}^{r} \in \mathscr{P}\left(U_{\alpha}\right)$. Because the generators can be uniquely expressed in terms of any other system of independent generators, there is for each pair $(\alpha, \beta)$ a $r \times r$ matrix

$$
G_{\alpha \beta}=\left(G_{\alpha \beta}^{j k}\right)_{1 \leqslant j, k \leqslant r}, \quad G_{\alpha \beta}^{j k} \in \mathscr{A}\left(U_{\alpha} \cap U_{\beta}\right),
$$

such that

$$
F_{\beta}^{k}=\sum_{1 \leqslant j \leqslant r} F_{\alpha}^{j} G_{\alpha \beta}^{j k} \quad \text { on } \quad U_{\alpha} \cap U_{\beta}
$$

In other words, we have a commutative diagram

$$
\begin{aligned}
& \mathscr{A}_{\mid U_{\alpha} \cap U_{\beta}}^{\oplus r} \xrightarrow{F_{\alpha}} \mathscr{S}_{\upharpoonright U_{\alpha} \cap U_{\beta}} \\
& G_{\alpha \beta} \uparrow \\
& \mathscr{A}_{\mid U_{\alpha} \cap U_{\beta}}^{\oplus r} \\
& \underset{F_{\beta}}{\longrightarrow} \mathscr{S}_{\uparrow U_{\alpha} \cap U_{\beta}}
\end{aligned}
$$

It follows easily from the equality $G_{\alpha \beta}=F_{\alpha}^{-1} \circ F_{\beta}$ that the transition matrices $G_{\alpha \beta}$ are invertible matrices satisfying the transition relation

$$
\begin{equation*}
G_{\alpha \gamma}=G_{\alpha \beta} G_{\beta \gamma} \quad \text { on } \quad U_{\alpha} \cap U_{\beta} \cap U_{\gamma} \tag{3.2}
\end{equation*}
$$

for all indices $\alpha, \beta, \gamma \in I$. In particular $G_{\alpha \alpha}=\operatorname{Id}$ on $U_{\alpha}$ and $G_{\alpha \beta}^{-1}=G_{\beta \alpha}$ on $U_{\alpha} \cap U_{\beta}$.
Conversely, if we are given a system of invertible $r \times r$ matrices $G_{\alpha \beta}$ with coefficients in $\mathcal{A}\left(U_{\alpha} \cap U_{\beta}\right)$ satisfying the transition relation (3.2), we can define a locally free sheaf $\mathscr{S}$ of rank $r$ over $\mathscr{A}$ by taking $\mathscr{S} \simeq \mathscr{A}^{\oplus r}$ over each $U_{\alpha}$, the identification over $U_{\alpha} \cap U_{\beta}$ being given by the isomorphism $G_{\alpha \beta}$. A section $H$ of $\mathscr{T}$ over an open set $\Omega \subset X$ can just be seen as a collection of sections $H_{\alpha}=\left(H_{\alpha}^{1}, \ldots, H_{\alpha}^{r}\right)$ of $\mathscr{A}^{\oplus r}\left(\Omega \cap U_{\alpha}\right)$ satisfying the transition relations $H_{\alpha}=G_{\alpha \beta} H_{\beta}$ over $\Omega \cap U_{\alpha} \cap U_{\beta}$.

The notion of locally free sheaf is closely related to another essential notion of differential geometry, namely the notion of vector bundle (resp. topological, differentiable, holomorphic ..., vector bundle). To describe the relation between these notions, we assume that the sheaf of rings $\mathscr{A}$ is a subsheaf of the sheaf $\mathscr{C}_{\mathbb{K}}$ of continous functions on $X$ with values in the field $\mathbb{K}=\mathbb{R}$ or $\mathbb{K}=\mathbb{C}$, containing the sheaf of locally constant functions $X \rightarrow \mathbb{K}$. Then, for each $x \in X$, there is an evaluation map

$$
\mathscr{A}_{x} \rightarrow \mathbb{K}, \quad w \mapsto w(x)
$$

whose kernel is a maximal ideal $\mathfrak{m}_{x}$ of $\mathscr{A}_{x}$, and $\mathscr{A}_{x} / \mathfrak{m}_{x}=\mathbb{K}$. Let $\mathscr{S}$ be a locally free sheaf of rank $r$ over $\mathscr{A}$. To each $x \in X$, we can associate a $\mathbb{K}$-vector space $E_{x}=\mathscr{S}_{x} / \mathfrak{m}_{x} \mathscr{S}_{x}$ : since $\mathscr{S}_{x} \simeq \mathscr{A}_{x}^{\oplus r}$, we have $E_{x} \simeq\left(\mathscr{A}_{x} / \mathfrak{m}_{x}\right)^{\oplus r}=\mathbb{K}^{r}$. The set $E=\coprod_{x \in X} E_{x}$ is equipped with a natural projection

$$
\pi: E \rightarrow X, \quad \xi \in E_{x} \mapsto \pi(\xi):=x
$$

and the fibers $E_{x}=\pi^{-1}(x)$ have a structure of $r$-dimensional $\mathbb{K}$-vector space: such a structure $E$ is called a $\mathbb{K}$-vector bundle of rank $r$ over $X$. Every section $s \in \mathscr{S}(U)$ gives rise to a section of $E$ over $U$ by setting $s(x)=s_{x} \bmod \mathfrak{m}_{x}$. We obtain a function (still denoted by the same symbol) $s: U \rightarrow E$ such that $s(x) \in E_{x}$ for every $x \in U$, i.e. $\pi \circ s=\operatorname{Id}_{U}$. It is clear that $\mathscr{S}(U)$ can be considered as a $\mathscr{A}(U)$-submodule of the $\mathbb{K}$-vector space of functions $U \rightarrow E$ mapping a point $x \in U$ to an element in the fiber $E_{x}$. Thus we get a subsheaf of the sheaf of $E$-valued sections, which is in a natural way a $\mathscr{A}$-module isomorphic to $\mathscr{P}$. This subsheaf will be denoted by $\mathscr{A}(E)$ and will be called the sheaf of $\mathscr{A}$-sections of $E$. If we are given a $\mathbb{K}$-vector bundle $E$ over $X$ and a subsheaf $\mathscr{S}=\mathscr{A}(E)$ of the sheaf of all sections of $E$ which is in a natural way a locally free $\mathscr{A}$-module of rank $r$, we say that $E$ (or more precisely the pair $(E, \mathcal{A}(E)$ )) is a $\mathscr{A}$-vector bundle of rank $r$ over $X$.
(3.3) Example. In case $\mathscr{A}=\mathscr{C}_{X, \mathbb{K}}$ is the sheaf of all $\mathbb{K}$-valued continuous functions on $X$, we say that $E$ is a topological vector bundle over $X$. When $X$ is a manifold and $\mathscr{A}=\mathscr{C}_{X, \mathrm{~K}}^{p}$, we say that $E$ is a $C^{p}$-differentiable vector bundle; finally, when $X$ is complex analytic and $\mathscr{A}=\mathscr{O}_{X}$, we say that $E$ is a holomorphic vector bundle.

Let us introduce still a little more notation. Since $\mathcal{A}(E)$ is a locally free sheaf of rank $r$ over any open set $U_{\alpha}$ in a suitable covering of $X$, a choice of generators ( $F_{\alpha}^{1}, \ldots, F_{\alpha}^{r}$ ) for $\mathscr{A}(E)_{\mid U_{\alpha}}$ yields corresponding generators $\left(e_{\alpha}^{1}(x), \ldots, e_{\alpha}^{r}(x)\right)$ of the fibers $E_{x}$ over $\mathbb{K}$. Such a system of generators is called a $\mathscr{A}$-admissible frame of $E$ over $U_{\alpha}$. There is a corresponding isomorphism

$$
\begin{equation*}
\theta_{\alpha}: E_{\mid U_{\alpha}}:=\pi^{-1}\left(U_{\alpha}\right) \longrightarrow U_{\alpha} \times \mathbb{K}^{r} \tag{3.4}
\end{equation*}
$$

which to each $\xi \in E_{x}$ associates the pair $\left(x,\left(\xi_{\alpha}^{1}, \ldots, \xi_{\alpha}^{r}\right)\right) \in U_{\alpha} \times \mathbb{K}^{r}$ composed of $x$ and of the components $\left(\xi_{\alpha}^{j}\right)_{1 \leqslant j \leqslant r}$ of $\xi$ in the basis $\left(e_{\alpha}^{1}(x), \ldots, e_{\alpha}^{r}(x)\right)$ of $E_{x}$. The bundle $E$ is said to be trivial if it is of the form $X \times \mathbb{K}^{r}$, which is the same as saying that $\mathscr{A}(E)=\mathscr{A}^{\oplus r}$. For this reason, the isomorphisms $\theta_{\alpha}$ are called trivializations of $E$ over $U_{\alpha}$. The corresponding transition automorphisms are

$$
\begin{aligned}
& \theta_{\alpha \beta}:=\theta_{\alpha} \circ \theta_{\beta}^{-1}:\left(U_{\alpha} \cap U_{\beta}\right) \times \mathbb{K}^{r} \longrightarrow\left(U_{\alpha} \cap U_{\beta}\right) \times \mathbb{K}^{r}, \\
& \theta_{\alpha \beta}(x, \xi)=\left(x, g_{\alpha \beta}(x) \cdot \xi\right), \quad(x, \xi) \in\left(U_{\alpha} \cap U_{\beta}\right) \times \mathbb{K}^{r},
\end{aligned}
$$

where $\left(g_{\alpha \beta}\right) \in \mathrm{GL}_{r}(\mathscr{A})\left(U_{\alpha} \cap U_{\beta}\right)$ are the transition matrices already described (except that they are just seen as matrices with coefficients in $\mathbb{K}$ rather than with coefficients in a sheaf). Conversely, if we are given a collection of matrices $g_{\alpha \beta}=\left(g_{\alpha \beta}^{j k}\right) \in \mathrm{GL}_{r}(\mathscr{A})\left(U_{\alpha} \cap\right.$ $U_{\beta}$ ) satisfying the transition relation

$$
g_{\alpha \gamma}=g_{\alpha \beta} g_{\beta \gamma} \quad \text { on } \quad U_{\alpha} \cap U_{\beta} \cap U_{\gamma},
$$

we can define a $\mathscr{A}$-vector bundle

$$
E=\left(\coprod_{\alpha \in I} U_{\alpha} \times \mathbb{K}^{r}\right) / \sim
$$

by gluing the charts $U_{\alpha} \times \mathbb{K}^{r}$ via the identification $\left(x_{\alpha}, \xi_{\alpha}\right) \sim\left(x_{\beta}, \xi_{\beta}\right)$ if and only if $x_{\alpha}=x_{\beta}=x \in U_{\alpha} \cap U_{\beta}$ and $\xi_{\alpha}=g_{\alpha \beta}(x) \cdot \xi_{\beta}$.
(3.5) Example. When $X$ is a real differentiable manifold, an interesting example of real vector bundle is the tangent bundle $T_{X}$; if $\tau_{\alpha}: U_{\alpha} \rightarrow \mathbb{R}^{n}$ is a collection of coordinate charts on $X$, then $\theta_{\alpha}=\pi \times d \tau_{\alpha}: T_{X \upharpoonright U_{\alpha}} \rightarrow U_{\alpha} \times \mathbb{R}^{m}$ define trivializations of $T_{X}$ and the transition matrices are given by $g_{\alpha \beta}(x)=d \tau_{\alpha \beta}\left(x^{\beta}\right)$ where $\tau_{\alpha \beta}=\tau_{\alpha} \circ \tau_{\beta}^{-1}$ and $x^{\beta}=\tau_{\beta}(x)$. The dual $T_{X}^{\star}$ of $T_{X}$ is called the cotangent bundle of $X$. If $X$ is complex analytic, then $T_{X}$ has the structure of a holomorphic vector bundle.

We now briefly discuss the concept of sheaf and bundle morphisms. If $\mathscr{I}$ and $\mathscr{S}^{\prime}$ are sheaves of $\mathscr{A}$-modules over a topological space $X$, then by a morphism $\varphi: \mathscr{S} \rightarrow \mathscr{S}^{\prime}$ we just mean a $\mathscr{A}$-linear sheaf morphism. If $\mathscr{S}=\mathscr{A}(E)$ and $\mathscr{S}^{\prime}=\mathscr{A}\left(E^{\prime}\right)$ are locally free sheaves, this is the same as a $\mathscr{A}$-linear bundle morphism, that is, a fiber preserving
$\mathbb{K}$-linear morphism $\varphi(x): E_{x} \rightarrow E_{x}^{\prime}$ such that the matrix representing $\varphi$ in any local $\mathscr{A}$-admissible frames of $E$ and $E^{\prime}$ has coefficients in $\mathscr{A}$.
(3.6) Proposition. Suppose that $\mathscr{A}$ is a sheaf of local rings, i.e. that a section of $\mathscr{A}$ is invertible in $\mathscr{A}$ if and only if it never takes the zero value in $\mathbb{K}$. Let $\varphi: \mathscr{S} \rightarrow \mathscr{S}^{\prime}$ be a $\mathscr{A}$-morphism of locally free $\mathscr{A}$-modules of rank $r, r^{\prime}$. If the rank of the $r^{\prime} \times r$ matrix $\varphi(x) \in M_{r^{\prime} r}(\mathbb{K})$ is constant for all $x \in X$, then $\operatorname{Ker} \varphi$ and $\operatorname{Im} \varphi$ are locally free subsheaves of $\mathscr{S}, \mathscr{S}^{\prime}$ respectively, and $\operatorname{Coker} \varphi=\mathscr{S}^{\prime} / \operatorname{Im} \varphi$ is locally free.

Proof. This is just a consequence of elementary linear algebra, once we know that non zero determinants with coefficients in $\mathscr{A}$ can be inverted.

Note that all three sheaves $\mathscr{C}_{X, \mathbb{K}}, \mathscr{C}_{X, \mathbb{K}}^{p}, \mathscr{O}_{X}$ are sheaves of local rings, so Prop. 3.6 applies to these cases. However, even if we work in the holomorphic category $\left(\mathscr{A}=\mathscr{O}_{X}\right)$, a difficulty immediately appears that the kernel or cokernel of an arbitrary morphism of locally free sheaves is in general not locally free.

## (3.7) Examples.

a) Take $X=\mathbb{C}$, let $\mathscr{S}=\mathscr{S}^{\prime}=\mathbb{O}$ be the trivial sheaf, and let $\varphi: \mathbb{O} \rightarrow \mathbb{O}$ be the morphism $u(z) \mapsto z u(z)$. It is immediately seen that $\varphi$ is injective as a sheaf morphism (© being an entire ring), and that Coker $\varphi$ is the skyscraper sheaf $\mathbb{C}_{0}$ of stalk $\mathbb{C}$ at $z=0$, having zero stalks at all other points $z \neq 0$. Thus $\operatorname{Coker} \varphi$ is not a locally free sheaf, although $\varphi$ is everywhere injective (note however that the corresponding morphism $\varphi: E \rightarrow E^{\prime},(z, \xi) \mapsto(z, z \xi)$ of trivial rank 1 vector bundles $E=E^{\prime}=\mathbb{C} \times \mathbb{C}$ is not injective on the zero fiber $E_{0}$ ).
b) Take $X=\mathbb{C}^{3}, \mathscr{S}=\mathscr{O}^{\oplus 3}, \mathscr{J}^{\prime}=\mathbb{O}$ and

$$
\varphi: \mathscr{O}^{\oplus 3} \rightarrow \mathbb{O}, \quad\left(u_{1}, u_{2}, u_{3}\right) \mapsto \sum_{1 \leqslant j \leqslant 3} z_{j} u_{j}\left(z_{1}, z_{2}, z_{3}\right) .
$$

Since $\varphi$ yields a surjective bundle morphism on $\mathbb{C}^{3} \backslash\{0\}$, one easily sees that $\operatorname{Ker} \varphi$ is locally free of rank 2 over $\mathbb{C}^{3} \backslash\{0\}$. However, by looking at the Taylor expansion of the $u_{j}$ 's at 0 , it is not difficult to check that $\operatorname{Ker} \varphi$ is the $\mathbb{O}$-submodule of $\mathbb{O}^{\oplus 3}$ generated by the three sections $\left(-z_{2}, z_{1}, 0\right),\left(-z_{3}, 0, z_{1}\right)$ and $\left(0, z_{3},-z_{2}\right)$, and that any two of these three sections cannot generate the 0 -stalk $(\operatorname{Ker} \varphi)_{0}$. Hence $\operatorname{Ker} \varphi$ is not locally free.

Since the category of locally free © $\mathbb{C}$-modules is not stable by taking kernels or cokernels, one is led to introduce a more general category which will be stable under these operations. This leads to the notion of coherent sheaves.

## $\S$ 3.2. Notion of Coherence

The notion of coherence again deals with sheaves of modules over a sheaf of rings. It is a semi-local property which says roughly that the sheaf of modules locally has a finite presentation in terms of generators and relations. We describe here some general properties of this notion, before concentrating ourselves on the case of coherent $\mathbb{O}_{X^{-}}$ modules.
(3.8) Definition. Let $\mathscr{A}$ be a sheaf of rings on a topological space $X$ and $\mathscr{S}$ a sheaf of modules over $\mathscr{A}$ (or briefly a $\mathscr{A}$-module). Then $\mathscr{S}$ is said to be locally finitely generated if for every point $x_{0} \in X$ one can find a neighborhood $\Omega$ and sections $F_{1}, \ldots, F_{q} \in \mathscr{S}(\Omega)$ such that for every $x \in \Omega$ the stalk $\mathscr{S}_{x}$ is generated by the germs $F_{1, x}, \ldots, F_{q, x}$ as an $A_{x}$-module.
(3.9) Lemma. Let $\mathscr{S}$ be a locally finitely generated sheaf of $\mathscr{A}$-modules on $X$ and $G_{1}, \ldots, G_{N}$ sections in $\mathscr{S}(U)$ such that $G_{1, x_{0}}, \ldots, G_{N, x_{0}}$ generate $\mathscr{S}_{x_{0}}$ at $x_{0} \in U$. Then $G_{1, x}, \ldots, G_{N, x}$ generate $\mathscr{S}_{x}$ for $x$ near $x_{0}$.

Proof. Take $F_{1}, \ldots, F_{q}$ as in Def. 3.8. As $G_{1}, \ldots, G_{N}$ generate $\mathscr{S}_{x_{0}}$, one can find a neighborhood $\Omega^{\prime} \subset \Omega$ of $x_{0}$ and $H_{j k} \in \mathscr{A}\left(\Omega^{\prime}\right)$ such that $F_{j}=\sum H_{j k} G_{k}$ on $\Omega^{\prime}$. Thus $G_{1, x}, \ldots, G_{N, x}$ generate $\mathscr{S}_{x}$ for all $x \in \Omega^{\prime}$.
$\S$ 3.2.1. Definition of Coherent Sheaves. If $U$ is an open subset of $X$, we denote by $\mathscr{S}_{\uparrow U}$ the restriction of $\mathscr{S}$ to $U$, i.e. the union of all stalks $\mathscr{S}_{x}$ for $x \in U$. If $F_{1}, \ldots, F_{q} \in \mathscr{S}(U)$, the kernel of the sheaf homomorphism $F: \mathscr{A}_{\uparrow U}^{\oplus q} \longrightarrow \mathscr{S}_{\mid U}$ defined by

$$
\begin{equation*}
\mathscr{A}_{x}^{\oplus q} \ni\left(g^{1}, \ldots, g^{q}\right) \longmapsto \sum_{1 \leqslant j \leqslant q} g^{j} F_{j, x} \in \mathscr{S}_{x}, \quad x \in U \tag{3.10}
\end{equation*}
$$

is a subsheaf $\mathscr{R}\left(F_{1}, \ldots, F_{q}\right)$ of $\mathscr{A}_{\lceil U}^{\oplus q}$, called the sheaf of relations between $F_{1}, \ldots, F_{q}$.
(3.11) Definition. A sheaf $\mathscr{I}$ of $\mathscr{A}$-modules on $X$ is said to be coherent if:
a) $\mathscr{S}$ is locally finitely generated ;
b) for any open subset $U$ of $X$ and any $F_{1}, \ldots, F_{q} \in \mathscr{S}(U)$, the sheaf of relations $\mathscr{R}\left(F_{1}, \ldots, F_{q}\right)$ is locally finitely generated.

Assumption a) means that every point $x \in X$ has a neighborhood $\Omega$ such that there is a surjective sheaf morphism $F: \mathscr{A}_{\uparrow \Omega}^{\oplus q} \longrightarrow \mathscr{S}_{\mid \Omega}$, and assumption b) implies that the kernel of $F$ is locally finitely generated. Thus, after shrinking $\Omega$, we see that $\mathscr{I}$ admits over $\Omega$ a finite presentation under the form of an exact sequence

$$
\begin{equation*}
\mathscr{A}_{\Gamma \Omega}^{\oplus p} \xrightarrow{G} \mathscr{A}_{\Gamma \Omega}^{\oplus q} \xrightarrow{F} \mathscr{S}_{\Gamma \Omega} \longrightarrow 0, \tag{3.12}
\end{equation*}
$$

where $G$ is given by a $q \times p$ matrix $\left(G_{j k}\right)$ of sections of $\mathscr{A}(\Omega)$ whose columns $\left(G_{j 1}\right), \ldots$, $\left(G_{j p}\right)$ are generators of $\mathscr{R}\left(F_{1}, \ldots, F_{q}\right)$.

It is clear that every locally finitely generated subsheaf of a coherent sheaf is coherent. From this we easily infer:
(3.13) Theorem. Let $\varphi: \mathscr{F} \longrightarrow \mathscr{G}$ be a $\mathscr{A}$-morphism of coherent sheaves. Then $\operatorname{Im} \varphi$ and $\operatorname{ker} \varphi$ are coherent.

Proof. Clearly $\operatorname{Im} \varphi$ is a locally finitely generated subsheaf of $\mathscr{G}$, so it is coherent. Let $x_{0} \in X$, let $F_{1}, \ldots, F_{q} \in \mathscr{F}(\Omega)$ be generators of $\mathscr{F}$ on a neighborhood $\Omega$ of $x_{0}$, and
$G_{1}, \ldots, G_{r} \in \mathscr{A}\left(\Omega^{\prime}\right)^{\oplus q}$ be generators of $\mathscr{R}\left(\varphi\left(F_{1}\right), \ldots, \varphi\left(F_{q}\right)\right)$ on a neighborhood $\Omega^{\prime} \subset \Omega$ of $x_{0}$. Then $\operatorname{ker} \varphi$ is generated over $\Omega^{\prime}$ by the sections

$$
H_{j}=\sum_{k=1}^{q} G_{j}^{k} F_{k} \in \mathscr{F}\left(\Omega^{\prime}\right), \quad 1 \leqslant j \leqslant r .
$$

(3.14) Theorem. Let $0 \longrightarrow \mathscr{F} \longrightarrow \mathscr{S} \longrightarrow \mathscr{G} \longrightarrow 0$ be an exact sequence of $\mathscr{A}$-modules. If two of the sheaves $\mathscr{F}, \mathscr{S}, \mathscr{\mathscr { G }}$ are coherent, then all three are coherent.

Proof. If $\mathscr{S}$ and $\mathscr{G}$ are coherent, then $\mathscr{F}=\operatorname{ker}(\mathscr{T} \rightarrow \mathscr{G})$ is coherent by Th. 3.13. If $\mathscr{S}$ and $\mathscr{F}$ are coherent, then $\mathscr{G}$ is locally finitely generated; to prove the coherence, let $G_{1}, \ldots, G_{q} \in \mathscr{G}(U)$ and $x_{0} \in U$. Then there is a neighborhood $\Omega$ of $x_{0}$ and sections $\tilde{G}_{1}, \ldots, \tilde{G}_{q} \in \mathscr{S}(\Omega)$ which are mapped to $G_{1}, \ldots, G_{q}$ on $\Omega$. After shrinking $\Omega$, we may assume also that $\mathscr{F}_{\uparrow \Omega}$ is generated by sections $F_{1}, \ldots, F_{p} \in \mathscr{F}(\Omega)$. Then $\mathscr{R}\left(G_{1}, \ldots, G_{q}\right)$ is the projection on the last $q$-components of $\mathscr{R}\left(F_{1}, \ldots, F_{p}, \tilde{G}_{1}, \ldots, \tilde{G}_{q}\right) \subset \mathscr{A}^{p+q}$, which is finitely generated near $x_{0}$ by the coherence of $\mathscr{S}$. Hence $\mathscr{R}\left(G_{1}, \ldots, G_{q}\right)$ is finitely generated near $x_{0}$ and $\mathscr{G}$ is coherent.

Finally, assume that $\mathscr{F}$ and $\mathscr{G}$ are coherent. Let $x_{0} \in X$ be any point, let $F_{1}, \ldots, F_{p} \in$ $\mathscr{F}(\Omega)$ and $G_{1}, \ldots, G_{q} \in \mathscr{G}(\Omega)$ be generators of $\mathscr{F}, \mathcal{G}$ on a neighborhood $\Omega$ of $x_{0}$. There is a neighborhood $\Omega^{\prime}$ of $x_{0}$ such that $G_{1}, \ldots, G_{q}$ admit liftings $\tilde{G}_{1}, \ldots, \tilde{G}_{q} \in \mathscr{S}\left(\Omega^{\prime}\right)$. Then $\left(F_{1}, \ldots, F_{q}, \tilde{G}_{1}, \ldots, \tilde{G}_{q}\right)$ generate $\mathscr{S}_{\mid \Omega^{\prime}}$, so $\mathscr{S}$ is locally finitely generated. Now, let $S_{1}, \ldots, S_{q}$ be arbitrary sections in $\mathscr{G}(U)$ and $\bar{S}_{1}, \ldots, \bar{S}_{q}$ their images in $\mathscr{G}(U)$. For any $x_{0} \in U$, the sheaf of relations $\mathscr{R}\left(\bar{S}_{1}, \ldots, \bar{S}_{q}\right)$ is generated by sections $P_{1}, \ldots, P_{s} \in \mathscr{A}(\Omega)^{\oplus q}$ on a small neighborhood $\Omega$ of $x_{0}$. Set $P_{j}=\left(P_{j}^{k}\right)_{1 \leqslant k \leqslant q}$. Then $H_{j}=P_{j}^{1} S_{1}+\ldots+P_{j}^{q} S_{q}$, $1 \leqslant j \leqslant s$, are mapped to 0 in $\mathscr{G}$ so they can be seen as sections of $\mathscr{F}$. The coherence of $\mathscr{F}$ shows that $\mathscr{R}\left(H_{1}, \ldots, H_{s}\right)$ has generators $Q_{1}, \ldots, Q_{t} \in \mathscr{A}\left(\Omega^{\prime}\right)^{s}$ on a small neighborhood $\Omega^{\prime} \subset \Omega$ of $x_{0}$. Then $\mathscr{R}\left(S_{1}, \ldots, S_{q}\right)$ is generated over $\Omega^{\prime}$ by $R_{j}=\sum Q_{j}^{k} P_{k} \in \mathscr{A}\left(\Omega^{\prime}\right)$, $1 \leqslant j \leqslant t$, and $\mathscr{T}$ is coherent.
(3.15) Corollary. If $\mathscr{F}$ and $\mathscr{G}$ are coherent subsheaves of a coherent analytic sheaf $\mathscr{S}$, the intersection $\mathscr{F} \cap \mathscr{G}$ is a coherent sheaf.

Proof. Indeed, the intersection sheaf $\mathscr{F} \cap \mathscr{G}$ is the kernel of the composite morphism $\mathscr{F} \longrightarrow \mathscr{S} \longrightarrow \mathscr{S} / \mathscr{G}$, and $\mathscr{S} / \mathscr{G}$ is coherent.
$\S$ 3.2.2. Coherent Sheaf of Rings. A sheaf of rings $\mathscr{A}$ is said to be coherent if it is coherent as a module over itself. By Def. 3.11, this means that for any open set $U \subset X$ and any sections $F_{j} \in \mathscr{A}(U)$, the sheaf of relations $\mathscr{R}\left(F_{1}, \ldots, F_{q}\right)$ is finitely generated. The above results then imply that all free modules $\mathscr{A}^{\oplus p}$ are coherent. As a consequence:
(3.16) Theorem. If $\mathscr{A}$ is a coherent sheaf of rings, any locally finitely generated subsheaf of $\mathscr{A}^{\oplus p}$ is coherent. In particular, if $\mathscr{S}$ is a coherent $\mathscr{A}$-module and $F_{1}, \ldots, F_{q} \in \mathscr{I}(U)$, the sheaf of relations $\mathscr{R}\left(F_{1}, \ldots, F_{q}\right) \subset \mathscr{A}^{\oplus q}$ is also coherent.

Let $\mathscr{S}$ be a coherent sheaf of modules over a coherent sheaf of ring $\mathscr{A}$. By an iteration of construction (3.12), we see that for every integer $m \geqslant 0$ and every point $x \in X$ there is a neighborhood $\Omega$ of $x$ on which there is an exact sequence of sheaves

$$
\begin{equation*}
\mathscr{A}_{\Gamma \Omega}^{\oplus p_{m}} \xrightarrow{F_{m}} \mathscr{A}_{\uparrow \Omega}^{\oplus p_{m-1}} \longrightarrow \cdots \longrightarrow \mathscr{A}_{\uparrow \Omega}^{\oplus p_{1}} \xrightarrow{F_{1}} \mathscr{A}_{\upharpoonright \Omega}^{\oplus p_{0}} \xrightarrow{F_{0}} \mathscr{S}_{\upharpoonright \Omega} \longrightarrow 0, \tag{3.17}
\end{equation*}
$$

where $F_{j}$ is given by a $p_{j-1} \times p_{j}$ matrix of sections in $\mathscr{A}(\Omega)$.

## $\S$ 3.3. Analytic Sheaves and the Oka Theorem

Many properties of holomorphic functions which will be considered in this book can be expressed in terms of sheaves. Among them, analytic sheaves play a central role. The Oka theorem [Oka 1950] asserting the coherence of the sheaf of holomorphic functions can be seen as a far-reaching deepening of the noetherian property seen in Sect. 1. The theory of analytic sheaves could not be presented without it.
(3.18) Definition. Let $M$ be a n-dimensional complex analytic manifold and let $\mathbb{O}_{M}$ be the sheaf of germs of analytic functions on M. An analytic sheaf over $M$ is by definition a sheaf $\mathscr{S}$ of modules over $\mathscr{O}_{M}$.
(3.19) Coherence theorem of Oka. The sheaf of rings $\mathscr{O}_{M}$ is coherent for any complex manifold $M$.

Let $F_{1}, \ldots, F_{q} \in \mathscr{O}(U)$. Since $\mathbb{O}_{M, x}$ is Noetherian, we already know that every stalk $\mathscr{R}\left(F_{1}, \ldots, F_{q}\right)_{x} \subset \mathscr{O}_{M, x}^{\oplus q}$ is finitely generated, but the important new fact expressed by the theorem is that the sheaf of relations is locally finitely generated, namely that the "same" generators can be chosen to generate each stalk in a neighborhood of a given point.

Proof. By induction on $n=\operatorname{dim}_{\mathbb{C}} M$. For $n=0$, the stalks $\mathbb{O}_{M, x}$ are equal to $\mathbb{C}$ and the result is trivial. Assume now that $n \geqslant 1$ and that the result has already been proved in dimension $n-1$. Let $U$ be an open set of $M$ and $F_{1}, \ldots, F_{q} \in \mathbb{O}_{M}(U)$. To show that $\mathscr{R}\left(F_{1}, \ldots, F_{q}\right)$ is locally finitely generated, we may assume that $U=\Delta=\Delta^{\prime} \times \Delta_{n}$ is a polydisk in $\mathbb{C}^{n}$ centered at $x_{0}=0$; after a change of coordinates and multiplication of $F_{1}, \ldots, F_{q}$ by invertible functions, we may also suppose that $F_{1}, \ldots, F_{q}$ are Weierstrass polynomials in $z_{n}$ with coefficients in $\mathbb{O}\left(\Delta^{\prime}\right)$. We need a lemma.
(3.20) Lemma. If $x=\left(x^{\prime}, x_{n}\right) \in \Delta$, the $\mathscr{O}_{\Delta, x}$-module $\mathscr{R}\left(F_{1}, \ldots, F_{q}\right)_{x}$ is generated by those of its elements whose components are germs of analytic polynomials in $\mathbb{O}_{\Delta^{\prime}, x^{\prime}}\left[z_{n}\right]$ with a degree in $z_{n}$ at most equal to $\mu$, the maximum of the degrees of $F_{1}, \ldots, F_{q}$.

Proof. Assume for example that $F_{q}$ is of the maximum degree $\mu$. By the Weierstrass preparation Th. 1.1 and Lemma 1.9 applied at $x$, we can write $F_{q, x}=f^{\prime} f^{\prime \prime}$ where $f^{\prime}, f^{\prime \prime} \in \mathbb{O}_{\Delta^{\prime}, x^{\prime}}\left[z_{n}\right], f^{\prime}$ is a Weierstrass polynomial in $z_{n}-x_{n}$ and $f^{\prime \prime}(x) \neq 0$. Let $\mu^{\prime}$ and $\mu^{\prime \prime}$ denote the degrees of $f^{\prime}$ and $f^{\prime \prime}$ with respect to $z_{n}$, so $\mu^{\prime}+\mu^{\prime \prime}=\mu$. Now, take $\left(g^{1}, \ldots, g^{q}\right) \in \mathscr{R}\left(F_{1}, \ldots, F_{q}\right)_{x}$. The Weierstrass division theorem gives

$$
g^{j}=F_{q, x} t^{j}+r^{j}, \quad j=1, \ldots, q-1,
$$

where $t^{j} \in \mathbb{O}_{\Delta, x}$ and $r^{j} \in \mathbb{O}_{\Delta^{\prime}, x^{\prime}}\left[z_{n}\right]$ is a polynomial of degree $<\mu^{\prime}$. For $j=q$, define $r^{q}=g^{q}+\sum_{1 \leqslant j \leqslant q-1} t^{j} F_{j, x}$. We can write

$$
\begin{equation*}
\left(g^{1}, \ldots, g^{q}\right)=\sum_{1 \leqslant j \leqslant q} t^{j}\left(0, \ldots, F_{q}, \ldots, 0,-F_{j}\right)_{x}+\left(r^{1}, \ldots, r^{q}\right) \tag{3.21}
\end{equation*}
$$

where $F_{q}$ is in the $j$-th position in the $q$-tuples of the summation. Since these $q$-tuples are in $\mathscr{R}\left(F_{1}, \ldots, F_{q}\right)_{x}$, we have $\left(r^{1}, \ldots, r^{q}\right) \in \mathscr{R}\left(F_{1}, \ldots, F_{q}\right)_{x}$, thus

$$
\sum_{1 \leqslant j \leqslant q-1} F_{j, x} r^{j}+f^{\prime} f^{\prime \prime} r^{q}=0 .
$$

As the sum is a polynomial in $z_{n}$ of degree $<\mu+\mu^{\prime}$, it follows from Lemma 1.9 that $f^{\prime \prime} r^{q}$ is a polynomial in $z_{n}$ of degree $<\mu$. Now we have

$$
\left(r^{1}, \ldots, r^{q}\right)=1 / f^{\prime \prime}\left(f^{\prime \prime} r^{1}, \ldots, f^{\prime \prime} r^{q}\right)
$$

where $f^{\prime \prime} r^{j}$ is of degree $<\mu^{\prime}+\mu^{\prime \prime}=\mu$. In combination with (3.21) this proves the lemma.

Proof of Theorem 3.19 (end).. If $g=\left(g^{1}, \ldots, g^{q}\right)$ is one of the polynomials of $\mathscr{R}\left(F_{1}, \ldots, F_{q}\right)_{x}$ described in Lemma 3.20, we can write

$$
g^{j}=\sum_{0 \leqslant k \leqslant \mu} u^{j k} z_{n}^{k}, \quad u^{j k} \in \mathbb{O}_{\Delta^{\prime}, x^{\prime}} .
$$

The condition for $\left(g^{1}, \ldots, g^{q}\right)$ to belong to $\mathscr{R}\left(F_{1}, \ldots, F_{q}\right)_{x}$ therefore consists of $2 \mu+1$ linear conditions for the germ $u=\left(u^{j k}\right)$ with coefficients in $\mathbb{O}\left(\Delta^{\prime}\right)$. By the induction hypothesis, $\mathscr{O}_{\Delta^{\prime}}$ is coherent and Th. 3.16 shows that the corresponding modules of relations are generated over $\mathbb{O}_{\Delta^{\prime}, x^{\prime}}$, for $x^{\prime}$ in a neighborhood $\Omega^{\prime}$ of 0 , by finitely many $(q \times \mu)$-tuples $U_{1}, \ldots, U_{N} \in \mathscr{O}\left(\Omega^{\prime}\right)^{q \mu}$. By Lemma 3.20, $\mathscr{R}\left(F_{1}, \ldots, F_{q}\right)_{x}$ is generated at every point $x \in \Omega=\Omega^{\prime} \times \Delta_{n}$ by the germs of the corresponding polynomials

$$
G_{l}(z)=\left(\sum_{0 \leqslant k \leqslant \mu} U_{l}^{j k}\left(z^{\prime}\right) z_{n}^{k}\right)_{1 \leqslant j \leqslant q}, \quad z \in \Omega, \quad 1 \leqslant l \leqslant N .
$$

(3.22) Strong Noetherian property. Let $\mathscr{F}$ be a coherent analytic sheaf on a complex manifold $M$ and let $\mathscr{F}_{1} \subset \mathscr{F}_{2} \subset \ldots$ be an increasing sequence of coherent subsheaves of $\mathscr{F}$. Then the sequence $\left(\mathscr{F}_{k}\right)$ is stationary on every compact subset of $M$.

Proof. Since $\mathscr{F}$ is locally a quotient of a free module $\mathscr{O}_{M}^{\oplus q}$, we can pull back the sequence to $\mathscr{O}_{M}^{\oplus q}$ and thus suppose $\mathscr{F}=\mathscr{O}_{M}$ (by easy reductions similar to those in the proof of Th. 3.14). Suppose $M$ connected and $\mathscr{F}_{k_{0}} \neq\{0\}$ for some index $k_{0}$ (otherwise, there is nothing to prove). By the analytic continuation theorem, we easily see that $\mathscr{F}_{k_{0}, x} \neq\{0\}$ for every $x \in M$. We can thus find a non zero Weierstrass polynomial $P \in \mathscr{F}_{k_{0}}(V)$, $\operatorname{deg}_{z_{n}} P\left(z^{\prime}, z_{n}\right)=\mu$, in a coordinate neighborhood $V=\Delta^{\prime} \times \Delta_{n}$ of any point $x \in M$. A division by $P$ shows that for $k \geqslant k_{0}$ and $x \in V$, all stalks $\mathscr{F}_{k, x}$ are generated by $P_{x}$ and by polynomials of degree $<\mu$ in $z_{n}$ with coefficients in $\mathscr{O}_{\Delta^{\prime}, x^{\prime}}$. Therefore, we can apply induction on $n$ to the coherent $\mathbb{O}_{\Delta^{\prime}}$-module

$$
\mathscr{F}^{\prime}=\mathscr{F} \cap\left\{Q \in \mathbb{O}_{\Delta^{\prime}}\left[z_{n}\right] ; \operatorname{deg} Q \leqslant \mu\right\}
$$

and its increasing sequence of coherent subsheaves $\mathscr{F}_{k}^{\prime}=\mathscr{F}_{k} \cap \mathscr{F}^{\prime}$.

## §4. Complex Analytic Sets. Local Properties

## § 4.1. Definition. Irreducible Components

A complex analytic set is a set which can be defined locally by finitely many holomorphic equations; such a set has in general singular points, because no assumption is
made on the differentials of the equations. We are interested both in the description of the singularities and in the study of algebraic properties of holomorphic functions on analytic sets. For a more detailed study than ours, we refer to H. Cartan's seminar [Cartan 1950], to the books of [Gunning-Rossi 1965], [Narasimhan 1966] or the recent book by [Grauert-Remmert 1984].
(4.1) Definition. Let $M$ be a complex analytic manifold. A subset $A \subset M$ is said to be an analytic subset of $M$ if $A$ is closed and if for every point $x_{0} \in A$ there exist a neighborhood $U$ of $x_{0}$ and holomorphic functions $g_{1}, \ldots, g_{n}$ in $\mathbb{O}(U)$ such that

$$
A \cap U=\left\{z \in U ; g_{1}(z)=\ldots=g_{N}(z)=0\right\}
$$

Then $g_{1}, \ldots, g_{N}$ are said to be (local) equations of $A$ in $U$.

It is easy to see that a finite union or intersection of analytic sets is analytic: if $\left(g_{j}^{\prime}\right)$, $\left(g_{k}^{\prime \prime}\right)$ are equations of $A^{\prime}, A^{\prime \prime}$ in the open set $U$, then the family of all products ( $g_{j}^{\prime} g_{k}^{\prime \prime}$ ) and the family $\left(g_{j}^{\prime}\right) \cup\left(g_{k}^{\prime \prime}\right)$ define equations of $A^{\prime} \cup A^{\prime \prime}$ and $A^{\prime} \cap A^{\prime \prime}$ respectively.
(4.2) Remark. Assume that $M$ is connected. The analytic continuation theorem shows that either $A=M$ or $A$ has no interior point. In the latter case, each piece $A \cap U=g^{-1}(0)$ is the set of points where the function $\log |g|^{2}=\log \left(\left|g_{1}\right|^{2}+\cdots+\left|g_{N}\right|^{2}\right) \in \operatorname{Psh}(U)$ takes the value $-\infty$, hence $A$ is pluripolar. In particular $M \backslash A$ is connected and every function $f \in \mathbb{O}(M \backslash A)$ that is locally bounded near $A$ can be extended to a function $\tilde{f} \in \mathbb{O}(M)$.

We focus now our attention on local properties of analytic sets. By definition, a germ of set at a point $x \in M$ is an equivalence class of elements in the power set $\mathscr{P}(M)$, with $A \sim B$ if there is an open neighborhood $V$ of $x$ such that $A \cap V=B \cap V$. The germ of a subset $A \subset M$ at $x$ will be denoted by $(A, x)$. We most often consider the case when $A \subset M$ is a analytic set in a neighborhood $U$ of $x$, and in this case we denote by $\mathscr{I}_{A, x}$ the ideal of germs $f \in \mathscr{O}_{M, x}$ which vanish on $(A, x)$. Conversely, if $\mathscr{G}=\left(g_{1}, \ldots, g_{N}\right)$ is an ideal of $\mathbb{O}_{M, x}$, we denote by $(V(\mathscr{I}), x)$ the germ at $x$ of the zero variety $V(\mathscr{Y})=\left\{z \in U ; g_{1}(z)=\ldots=g_{N}(z)=0\right\}$, where $U$ is a neighborhood of $x$ such that $g_{j} \in \mathscr{G}(U)$. It is easy to check that the germ $(V(c J), x)$ does not depend on the choice of generators. Moreover, it is clear that

$$
\text { for every ideal } \mathscr{f} \text { in the ring } \mathscr{O}_{M, x}, \quad \mathscr{I}_{V(\mathscr{J}), x} \supset \mathscr{f},
$$

for every germ of analytic set $(A, x), \quad\left(V\left(\mathscr{I}_{A, x}\right), x\right)=(A, x)$.
(4.4) Definition. $A$ germ $(A, x)$ is said to be irreducible if it has no decomposition $(A, x)=\left(A_{1}, x\right) \cup\left(A_{2}, x\right)$ with analytic sets $\left(A_{j}, x\right) \neq(A, x), j=1,2$.
(4.5) Proposition. $A \operatorname{germ}(A, x)$ is irreducible if and only if $\mathscr{I}_{A, x}$ is a prime ideal of the ring $\mathbb{O}_{M, x}$.

Proof. Let us recall that an ideal $\mathscr{g}$ is said to be prime if $f g \in \mathscr{G}$ implies $f \in \mathscr{G}$ or $g \in \mathscr{g}$. Assume that $(A, x)$ is irreducible and that $f g \in \mathscr{I}_{A, x}$. As we can write $(A, x)=\left(A_{1}, x\right) \cup\left(A_{2}, x\right)$ with $A_{1}=A \cap f^{-1}(0)$ and $A_{2}=A \cap g^{-1}(0)$, we must have
for example $\left(A_{1}, x\right)=(A, x)$; thus $f \in \mathscr{I}_{A, x}$ and $\mathscr{I}_{A, x}$ is prime. Conversely, if $(A, x)=$ $\left(A_{1}, x\right) \cup\left(A_{2}, x\right)$ with $\left(A_{j}, x\right) \neq(A, x)$, there exist $f \in \mathscr{I}_{A_{1}, x}, g \in \mathscr{I}_{A_{2}, x}$ such that $f, g \notin \mathscr{I}_{A, x}$. However $f g \in \mathscr{I}_{A, x}$, thus $\mathscr{I}_{A, x}$ is not prime.
(4.6) Theorem. Every decreasing sequence of germs of analytic sets $\left(A_{k}, x\right)$ is stationary.

Proof. In fact, the corresponding sequence of ideals $\mathscr{g}_{k}=\mathscr{I}_{A_{k}, x}$ is increasing, thus $\mathscr{F}_{k}=\mathscr{F}_{k_{0}}$ for $k \geqslant k_{0}$ large enough by the Noetherian property of $\mathbb{O}_{M, x}$. Hence $\left(A_{k}, x\right)=\left(V\left(\mathscr{F}_{k}\right), x\right)$ is constant for $k \geqslant k_{0}$. This result has the following straightforward consequence:
(4.7) Theorem. Every analytic germ $(A, x)$ has a finite decomposition

$$
(A, x)=\bigcup_{1 \leqslant k \leqslant N}\left(A_{k}, x\right)
$$

where the germs $\left(A_{j}, x\right)$ are irreducible and $\left(A_{j}, x\right) \not \subset\left(A_{k}, x\right)$ for $j \neq k$. The decomposition is unique apart from the ordering.

Proof. If $(A, x)$ can be split in several components, we split repeatedly each component as long as one of them is reducible. The process must stop by Th. 4.6, whence the existence. For the uniqueness, assume that $(A, x)=\bigcup\left(A_{l}^{\prime}, x\right), 1 \leqslant l \leqslant N^{\prime}$, is another decomposition. Since $\left(A_{k}, x\right)=\bigcup_{l}\left(A_{k} \cap A_{l}^{\prime}, x\right)$, we must have $\left(A_{k}, x\right)=\left(A_{k} \cap A_{l}^{\prime}, x\right)$ for some $l=l(k)$, i.e. $\left(A_{k}, x\right) \subset\left(A_{l(k)}^{\prime}, x\right)$, and likewise $\left(A_{l(k)}^{\prime}, x\right) \subset\left(A_{j}, x\right)$ for some $j$. Hence $j=k$ and $\left(A_{l(k)}^{\prime}, x\right)=\left(A_{k}, x\right)$.

## § 4.2. Local Structure of a Germ of Analytic Set

We are going to describe the local structure of a germ, both from the holomorphic and topological points of view. By the above decomposition theorem, we may restrict ourselves to the case of irreducible germs Let $\mathscr{g}$ be a prime ideal of $\mathbb{O}_{n}=\mathbb{O}_{\mathbb{C}^{n}, 0}$ and let $A=V(\mathscr{G})$ be its zero variety. We set $\mathscr{g}_{k}=\mathscr{g} \cap \mathbb{C}\left\{z_{1}, \ldots, z_{k}\right\}$ for each $k=0,1, \ldots, n$.
(4.8) Proposition. There exist an integer d, a basis $\left(e_{1}, \ldots, e_{n}\right)$ of $\mathbb{C}^{n}$ and associated coordinates $\left(z_{1}, \ldots, z_{n}\right)$ with the following properties: $\mathscr{H}_{d}=\{0\}$ and for every integer $k=d+1, \ldots, n$ there is a Weierstrass polynomial $P_{k} \in \mathscr{g}_{k}$ of the form

$$
\begin{equation*}
P_{k}\left(z^{\prime}, z_{k}\right)=z_{k}^{s_{k}}+\sum_{1 \leqslant j \leqslant s_{k}} a_{j, k}\left(z^{\prime}\right) z_{k}^{s_{k}-j}, \quad a_{j, k}\left(z^{\prime}\right) \in \mathscr{O}_{k-1}, \tag{4.9}
\end{equation*}
$$

where $a_{j, k}\left(z^{\prime}\right)=O\left(\left|z^{\prime}\right|^{j}\right)$. Moreover, the basis $\left(e_{1}, \ldots, e_{n}\right)$ can be chosen arbitrarily close to any preassigned basis $\left(e_{1}^{0}, \ldots, e_{n}^{0}\right)$.

Proof. By induction on $n$. If $\mathscr{f}=\mathscr{I}_{n}=\{0\}$, then $d=n$ and there is nothing to prove. Otherwise, select a non zero element $g_{n} \in \mathscr{g}$ and a vector $e_{n}$ such that $\mathbb{C} \ni w \longmapsto$ $g_{n}\left(w e_{n}\right)$ has minimum vanishing order $s_{n}$. This choice excludes at most the algebraic set $g_{n}^{\left(s_{n}\right)}(v)=0$, so $e_{n}$ can be taken arbitrarily close to $e_{n}^{0}$. Let $\left(\tilde{z}_{1}, \ldots, \tilde{z}_{n-1}, z_{n}\right)$ be the coordinates associated to the basis $\left(e_{1}^{0}, \ldots, e_{n-1}^{0}, e_{n}\right)$. After multiplication by an invertible element, we may assume that $g_{n}$ is a Weierstrass polynomial

$$
P_{n}\left(\tilde{z}, z_{n}\right)=z_{n}^{s_{n}}+\sum_{1 \leqslant j \leqslant s_{n}} a_{j, n}(\tilde{z}) z_{n}^{s_{n}-j}, \quad a_{j, n} \in \mathbb{O}_{n-1},
$$

and $a_{j, n}(\tilde{z})=O\left(|\tilde{z}|^{j}\right)$ by Remark 2.2. If $\mathscr{g}_{n-1}=\mathscr{G} \cap \mathbb{C}\{\tilde{z}\}=\{0\}$ then $d=n-1$ and the construction is finished. Otherwise we apply the induction hypothesis to the ideal $\mathscr{F}_{n-1} \subset \mathscr{O}_{n-1}$ in order to find a new basis $\left(e_{1}, \ldots, e_{n-1}\right)$ of $\operatorname{Vect}\left(e_{1}^{0}, \ldots, e_{n-1}^{0}\right)$, associated coordinates $\left(z_{1}, \ldots, z_{n-1}\right)$ and Weierstrass polynomials $P_{k} \in \mathscr{F}_{k}, d+1 \leqslant k \leqslant n-1$, as stated in the lemma.
(4.10) Lemma. If $w \in \mathbb{C}$ is a root of $w^{d}+a_{1} w^{d-1}+\cdots+a_{d}=0, a_{j} \in \mathbb{C}$, then $|w| \leqslant 2 \max \left|a_{j}\right|^{1 / j}$.

Proof. Otherwise $|w|>2\left|a_{j}\right|^{1 / j}$ for all $j=1, \ldots, d$ and the given equation $-1=a_{1} / w+$ $\cdots+a_{d} / w^{d}$ implies $1 \leqslant 2^{-1}+\cdots+2^{-d}$, a contradiction.
(4.11) Corollary. Set $z^{\prime}=\left(z_{1}, \ldots, z_{d}\right), z^{\prime \prime}=\left(z_{d+1}, \ldots, z_{n}\right)$, and let $\Delta^{\prime}$ in $\mathbb{C}^{d}$, $\Delta^{\prime \prime}$ in $\mathbb{C}^{n-d}$ be polydisks of center 0 and radii $r^{\prime}, r^{\prime \prime}>0$. Then the germ $(A, 0)$ is contained in a cone $\left|z^{\prime \prime}\right| \leqslant C\left|z^{\prime}\right|, C=$ constant, and the restriction of the projection map $\mathbb{C}^{n} \longrightarrow \mathbb{C}^{d}$, $\left(z^{\prime}, z^{\prime \prime}\right) \longmapsto z^{\prime} \quad:$

$$
\pi: A \cap\left(\Delta^{\prime} \times \Delta^{\prime \prime}\right) \longrightarrow \Delta^{\prime}
$$

is proper if $r^{\prime \prime}$ is small enough and $r^{\prime} \leqslant r^{\prime \prime} / C$.
Proof. The polynomials $P_{k}\left(z_{1}, \ldots, z_{k-1} ; z_{k}\right)$ vanish on ( $A, 0$ ). By Lemma 4.10 and (4.9), every point $z \in A$ sufficiently close to 0 satisfies

$$
\left|z_{k}\right| \leqslant C_{k}\left(\left|z_{1}\right|+\cdots+\left|z_{k-1}\right|\right), \quad d+1 \leqslant k \leqslant n
$$

thus $\left|z^{\prime \prime}\right| \leqslant C\left|z^{\prime}\right|$ and the Corollary follows.
Since $\mathscr{G}_{d}=\{0\}$, we have an injective ring morphism

$$
\begin{equation*}
\mathfrak{O}_{d}=\mathbb{C}\left\{z_{1}, \ldots, z_{d}\right\} \hookrightarrow \mathscr{O}_{n} / \mathscr{G} . \tag{4.12}
\end{equation*}
$$

(4.13) Proposition. $\mathscr{O}_{n} / \mathscr{g}$ is a finite integral extension of $\mathscr{O}_{d}$.

Proof. Let $f \in \mathcal{O}_{n}$. A division by $P_{n}$ yields $f=P_{n} q_{n}+R_{n}$ with a remainder $R_{n} \in$ $\mathfrak{O}_{n-1}\left[z_{n}\right], \operatorname{deg}_{z_{n}} R_{n}<s_{n}$. Further divisions of the coefficients of $R_{n}$ by $P_{n-1}, P_{n-2}$ etc ... yield

$$
R_{k+1}=P_{k} q_{k}+R_{k}, \quad R_{k} \in \mathscr{O}_{k}\left[z_{k+1}, \ldots, z_{n}\right],
$$

where $\operatorname{deg}_{z_{j}} R_{k}<s_{j}$ for $j>k$. Hence

$$
\begin{equation*}
f=R_{d}+\sum_{d+1 \leqslant k \leqslant n} P_{k} q_{k}=R_{d} \quad \bmod \quad\left(P_{d+1}, \ldots, P_{n}\right) \subset \mathscr{g} \tag{4.14}
\end{equation*}
$$

and $\mathscr{O}_{n} / \mathscr{g}$ is finitely generated as an $\mathscr{O}_{d}$-module by the family of monomials $z_{d+1}^{\alpha_{d+1}} \ldots z_{n}^{\alpha_{n}}$ with $\alpha_{j}<s_{j}$.

As $\mathscr{g}$ is prime, $\mathscr{O}_{n} / \mathscr{g}$ is an entire ring. We denote by $\tilde{f}$ the class of any germ $f \in \mathscr{O}_{n}$ in $\mathscr{O}_{n} / \mathscr{F}$, by $\mathscr{M}_{A}$ and $\mathscr{M}_{d}$ the quotient fields of $\mathscr{O}_{n} / \mathscr{\mathscr { L }}$ and $\mathscr{O}_{d}$ respectively. Then $\mathcal{M}_{A}=\mathscr{M}_{d}\left[\tilde{z}_{d+1}, \ldots, \tilde{z}_{n}\right]$ is a finite algebraic extension of $\mathscr{M}_{d}$. Let $q=\left[\mathcal{M}_{A}: \mathcal{M}_{d}\right]$ be its degree and let $\sigma_{1}, \ldots, \sigma_{q}$ be the embeddings of $\mathcal{M}_{A}$ over $\mathscr{M}_{d}$ in an algebraic closure $\overline{\mathcal{M}}_{A}$. Let us recall that a factorial ring is integrally closed in its quotient field ([Lang 1965],

Chapter IX). Hence every element of $\mathscr{M}_{d}$ which is integral over $\mathscr{O}_{d}$ lies in fact in $\mathscr{O}_{d}$. By the primitive element theorem, there exists a linear form $u\left(z^{\prime \prime}\right)=c_{d+1} z_{d+1}+\cdots+c_{n} z_{n}$, $c_{k} \in \mathbb{C}$, such that $\mathscr{M}_{A}=\mathscr{M}_{d}[\tilde{u}]$; in fact, $u$ is of degree $q$ if and only if $\sigma_{1} \tilde{u}, \ldots, \sigma_{q} \tilde{u}$ are all distinct, and this excludes at most a finite number of vector subspaces in the space $\mathbb{C}^{n-d}$ of coefficients $\left(c_{d+1}, \ldots, c_{n}\right)$. As $\tilde{u} \in \mathscr{O}_{n} / \mathscr{F}$ is integral over the integrally closed ring $\mathbb{O}_{d}$, the unitary irreducible polynomial $W_{u}$ of $\tilde{u}$ over $\mathcal{M}_{d}$ has coefficients in $\mathscr{O}_{d}$ :

$$
W_{u}\left(z^{\prime} ; T\right)=T^{q}+\sum_{1 \leqslant j \leqslant q} a_{j}\left(z_{1}, \ldots, z_{d}\right) T^{q-j}, \quad a_{j} \in \mathbb{O}_{d}
$$

$W_{u}$ must be a Weierstrass polynomial, otherwise there would exist a factorization $W_{u}=$ $W^{\prime} Q$ in $\mathscr{O}_{d}[T]$ with a Weierstrass polynomial $W^{\prime}$ of degree $\operatorname{deg} W^{\prime}<q=\operatorname{deg} \tilde{u}$ and $Q(0) \neq 0$, hence $W^{\prime}(\tilde{u})=0$, a contradiction. In the same way, we see that $\tilde{z}_{d+1}, \ldots, \tilde{z}_{n}$ have irreducible equations $W_{k}\left(z^{\prime} ; \tilde{z}_{k}\right)=0$ where $W_{k} \in \mathbb{O}_{d}[T]$ is a Weierstrass polynomial of degree $=\operatorname{deg} \tilde{z}_{k} \leqslant q, d+1 \leqslant k \leqslant n$.
(4.15) Lemma. Let $\delta\left(z^{\prime}\right) \in \mathfrak{O}_{d}$ be the discriminant of $W_{u}\left(z^{\prime} ; T\right)$. For every element $g$ of $\mathscr{M}_{A}$ which is integral over $\mathfrak{O}_{d}$ (or equivalently over $\mathscr{O}_{n} / \mathscr{g}$ ) we have $\delta g \in \mathbb{O}_{d}[\tilde{u}]$.

Proof. We have $\delta\left(z^{\prime}\right)=\prod_{j<k}\left(\sigma_{k} \tilde{u}-\sigma_{j} \tilde{u}\right)^{2} \not \equiv 0$, and $g \in \mathcal{M}_{A}=\mathcal{M}_{d}[\tilde{u}]$ can be written

$$
g=\sum_{0 \leqslant j \leqslant q-1} b_{j} \tilde{u}^{j}, \quad b_{j} \in \mathcal{M}_{d}
$$

where $b_{0}, \ldots, b_{d-1}$ are the solutions of the linear system $\sigma_{k} g=\sum b_{j}\left(\sigma_{k} \tilde{u}\right)^{j}$; the determinant (of Van der Monde type) is $\delta^{1 / 2}$. It follows that $\delta b_{j} \in \mathcal{M}_{d}$ are polynomials in $\sigma_{k} g$ and $\sigma_{k} \tilde{u}$, thus $\delta b_{j}$ is integral over $\mathscr{O}_{d}$. As $\mathscr{O}_{d}$ is integrally closed, we must have $\delta b_{j} \in \mathscr{O}_{d}$, hence $\delta g \in \mathscr{O}_{d}[\tilde{u}]$.

In particular, there exist unique polynomials $B_{d+1}, \ldots, B_{n} \in \mathscr{O}_{d}[T]$ with $\operatorname{deg} B_{k} \leqslant$ $q-1$, such that

$$
\begin{equation*}
\delta\left(z^{\prime}\right) z_{k}=B_{k}\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right) \quad(\bmod \mathscr{y}) \tag{4.16}
\end{equation*}
$$

Then we have

$$
\begin{equation*}
\delta\left(z^{\prime}\right)^{q} W_{k}\left(z^{\prime} ; B_{k}\left(z^{\prime} ; T\right) / \delta\left(z^{\prime}\right)\right) \in \text { ideal } W_{u}\left(z^{\prime} ; T\right) \mathscr{O}_{d}[T] ; \tag{4.17}
\end{equation*}
$$

indeed, the left-hand side is a polynomial in $\mathscr{O}_{d}[T]$ and admits $T=\tilde{u}$ as a root in $\mathscr{O}_{n} / \mathscr{G}$ since $B_{k}\left(z^{\prime} ; \tilde{u}\right) / \delta\left(z^{\prime}\right)=\tilde{z}_{k}$ and $W_{k}\left(z^{\prime} ; \tilde{z}_{k}\right)=0$.
(4.18) Lemma. Consider the ideal

$$
\mathscr{G}=\left(W_{u}\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right), \delta\left(z^{\prime}\right) z_{k}-B_{k}\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right)\right) \subset \mathscr{y}
$$

and set $m=\max \{q,(n-d)(q-1)\}$. For every germ $f \in \mathscr{O}_{n}$, there exists a unique polynomial $R \in \mathscr{O}_{d}[T], \operatorname{deg}_{T} R \leqslant q-1$, such that

$$
\delta\left(z^{\prime}\right)^{m} f(z)=R\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right) \quad(\bmod \mathscr{G}) .
$$

Moreover $f \in \mathscr{G}$ implies $R=0$, hence $\delta^{m} \mathscr{G} \subset \mathscr{G}$.
Proof. By (4.17) and a substitution of $z_{k}$, we find $\delta\left(z^{\prime}\right)^{q} W_{k}\left(z^{\prime} ; z_{k}\right) \in \mathscr{G}$. The analogue of formula (4.14) with $W_{k}$ in place of $P_{k}$ yields

$$
f=R_{d}+\sum_{d+1 \leqslant k \leqslant n} W_{k} q_{k}, \quad R_{d} \in \mathscr{O}_{d}\left[z_{d+1}, \ldots, z_{n}\right],
$$

with $\operatorname{deg}_{z_{k}} R_{d}<\operatorname{deg} W_{k} \leqslant q$, thus $\delta^{m} f=\delta^{m} R_{d} \bmod \mathscr{G}$. We may therefore replace $f$ by $R_{d}$ and assume that $f \in \mathbb{O}_{d}\left[z_{d+1}, \ldots, z_{n}\right]$ is a polynomial of total degree $\leqslant(n-d)(q-1) \leqslant$ $m$. A substitution of $z_{k}$ by $B_{k}\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right) / \delta\left(z^{\prime}\right)$ yields $G \in \mathscr{O}_{d}[T]$ such that

$$
\delta\left(z^{\prime}\right)^{m} f(z)=G\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right) \quad \bmod \quad\left(\delta\left(z^{\prime}\right) z_{k}-B_{k}\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right)\right) .
$$

Finally, a division $G=W_{u} Q+R$ gives the required polynomial $R \in \mathscr{O}_{d}[T]$. The last statement is clear: if $f \in \mathscr{G}$ satisfies $\delta^{m}\left(z^{\prime}\right) f(z)=R\left(z ; u\left(z^{\prime \prime}\right)\right) \bmod \mathscr{G}$ for $\operatorname{deg}_{T} R<q$, then $R\left(z^{\prime} ; \tilde{u}\right)=0$, and as $\tilde{u} \in \mathscr{O}_{n} / \mathscr{G}$ is of degree $q$, we must have $R=0$. The uniqueness of $R$ is proved similarly.
(4.19) Local parametrization theorem. Let $\mathscr{g}$ be a prime ideal of $\mathbb{O}_{n}$ and let $A=$ $V(\mathscr{y})$. Assume that the coordinates

$$
\left(z^{\prime} ; z^{\prime \prime}\right)=\left(z_{1}, \ldots, z_{d} ; z_{d+1}, \ldots, z_{n}\right)
$$

are chosen as above. Then the ring $\mathbb{O}_{n} / \mathscr{G}$ is a finite integral extension of $\mathbb{O}_{d}$; let $q$ be the degree of the extension and let $\delta\left(z^{\prime}\right) \in \mathbb{O}_{d}$ be the discriminant of the irreducible polynomial of a primitive element $u\left(z^{\prime \prime}\right)=\sum_{k>d} c_{k} z_{k}$. If $\Delta^{\prime}, \Delta^{\prime \prime}$ are polydisks of sufficiently small radii $r^{\prime}, r^{\prime \prime}$ and if $r^{\prime} \leqslant r^{\prime \prime} / C$ with $C$ large, the projection map $\pi: A \cap\left(\Delta^{\prime} \times \Delta^{\prime \prime}\right) \longrightarrow \Delta^{\prime}$ is a ramified covering with $q$ sheets, whose ramification locus is contained in $S=\left\{z^{\prime} \in\right.$ $\left.\Delta^{\prime} ; \delta\left(z^{\prime}\right)=0\right\}$. This means that:
a) the open subset $A_{S}=A \cap\left(\left(\Delta^{\prime} \backslash S\right) \times \Delta^{\prime \prime}\right)$ is a smooth d-dimensional manifold, dense in $A \cap\left(\Delta^{\prime} \times \Delta^{\prime \prime}\right)$;
b) $\pi: A_{S} \longrightarrow \Delta^{\prime} \backslash S$ is a covering ;
c) the fibers $\pi^{-1}\left(z^{\prime}\right)$ have exactly $q$ elements if $z^{\prime} \notin S$ and at most $q$ if $z^{\prime} \in S$.

Moreover, $A_{S}$ is a connected covering of $\Delta^{\prime} \backslash S$, and $A \cap\left(\Delta^{\prime} \times \Delta^{\prime \prime}\right)$ is contained in a cone $\left|z^{\prime \prime}\right| \leqslant C\left|z^{\prime}\right|$ (see Fig. 1).


Fig. II-1 Ramified covering from $A$ to $\Delta^{\prime} \subset \mathbb{C}^{p}$.

Proof. After a linear change in the coordinates $z_{d+1}, \ldots, z_{n}$, we may assume $u\left(z^{\prime \prime}\right)=z_{d+1}$, so $W_{u}=W_{d+1}$ and $B_{d+1}\left(z^{\prime} ; T\right)=\delta\left(z^{\prime}\right) T$. By Lemma 4.18, we have

$$
\mathscr{G}=\left(W_{d+1}\left(z^{\prime}, z_{d+1}\right), \delta\left(z^{\prime}\right) z_{k}-B_{k}\left(z^{\prime}, z_{d+1}\right)\right)_{k \geqslant d+2} \subset \mathscr{y}, \quad \delta^{m} \mathscr{G} \subset \mathscr{G} .
$$

We can thus find a polydisk $\Delta=\Delta^{\prime} \times \Delta^{\prime \prime}$ of sufficiently small radii $r^{\prime}, r^{\prime \prime}$ such that $V(\mathscr{F}) \subset V(\mathscr{G}) \subset V\left(\delta^{m} \mathscr{F}\right)$ in $\Delta$. As $V(\mathscr{G})=A$ and $V(\delta) \cap \Delta=S \times \Delta^{\prime \prime}$, this implies

$$
A \cap \Delta \subset V(\mathscr{G}) \cap \Delta \subset(A \cap \Delta) \cup\left(S \times \Delta^{\prime \prime}\right)
$$

In particular, the set $A_{S}=A \cap\left(\left(\Delta^{\prime} \backslash S\right) \times \Delta^{\prime \prime}\right)$ lying above $\Delta^{\prime} \backslash S$ coincides with $V(\mathscr{G}) \cap\left(\left(\Delta^{\prime} \backslash S\right) \times \Delta^{\prime \prime}\right)$, which is the set of points $z \in \Delta$ parametrized by the equations

$$
\left\{\begin{array}{l}
\delta\left(z^{\prime}\right) \neq 0, \quad W_{d+1}\left(z^{\prime}, z_{d+1}\right)=0,  \tag{4.20}\\
z_{k}=B_{k}\left(z^{\prime}, z_{d+1}\right) / \delta\left(z^{\prime}\right), \quad d+2 \leqslant k \leqslant n .
\end{array}\right.
$$

As $\delta\left(z^{\prime}\right)$ is the resultant of $W_{d+1}$ and $\partial W_{d+1} / \partial T$, we have

$$
\partial W_{d+1} / \partial T\left(z^{\prime}, z_{d+1}\right) \neq 0 \quad \text { on } \quad A_{S} .
$$

The implicit function theorem shows that $z_{d+1}$ is locally a holomorphic function of $z^{\prime}$ on $A_{S}$, and the same is true for $z_{k}=B_{k}\left(z^{\prime}, z_{d+1}\right) / \delta\left(z^{\prime}\right), k \geqslant d+2$. Hence $A_{S}$ is a smooth manifold, and for $r^{\prime} \leqslant r^{\prime \prime} / C$ small, the projection map $\pi: A_{S} \longrightarrow \Delta^{\prime} \backslash S$ is a proper local diffeomorphism; by (4.20) the fibers $\pi^{-1}\left(z^{\prime}\right)$ have at most $q$ points corresponding to some of the $q$ roots $w$ of $W_{d+1}\left(z^{\prime} ; w\right)=0$. Since $\Delta^{\prime} \backslash S$ is connected (Remark 4.2), either $A_{S}=\emptyset$ or the map $\pi$ is a covering of constant sheet number $q^{\prime} \leqslant q$. However, if $w$ is a root of $W_{d+1}\left(z^{\prime}, w\right)=0$ with $z^{\prime} \in \Delta^{\prime} \backslash S$ and if we set $z_{d+1}=w, z_{k}=B_{k}\left(z^{\prime}, w\right) / \delta\left(z^{\prime}\right)$, $k \geqslant d+2$, relation (4.17) shows that $W_{k}\left(z^{\prime}, z_{k}\right)=0$, in particular $\left|z_{k}\right|=O\left(\left|z^{\prime}\right|^{1 / q}\right)$ by Lemma 4.10. For $z^{\prime}$ small enough, the $q$ points $z=\left(z^{\prime}, z^{\prime \prime}\right)$ defined in this way lie in $\Delta$, thus $q^{\prime}=q$. Property b) and the first parts of a) and c) follow. Now, we need the following lemma.
(4.21) Lemma. If $\mathscr{G} \subset \mathscr{O}_{n}$ is prime and $A=V(\mathscr{G})$, then $\mathscr{I}_{A, 0}=\mathscr{G}$.
I.t is obvious that $\mathscr{I}_{A, 0} \supset \mathscr{G}$. Now, for any $f \in \mathscr{I}_{A, 0}$, Prop. 4.13 implies that $\tilde{f}$ satisfies in $\mathfrak{O}_{n} / \mathscr{I}$ an irreducible equation

$$
f^{r}+b_{1}\left(z^{\prime}\right) f^{r-1}+\cdots+b_{r}\left(z^{\prime}\right)=0 \quad(\bmod \mathscr{y}) .
$$

Then $b_{r}\left(z^{\prime}\right)$ vanishes on $(A, 0)$ and the first part of c) gives $b_{r}=0$ on $\Delta^{\prime} \backslash S$. Hence $\tilde{b}_{r}=0$ and the irreducibility of the equation of $\tilde{f}$ implies $r=1$, so $f \in \mathscr{G}$, as desired.

Proof of Theorem 4.19 (end).. It only remains to prove that $A_{S}$ is connected and dense in $A \cap \Delta$ and that the fibers $\pi^{-1}\left(z^{\prime}\right), z^{\prime} \in S$, have at most $q$ elements. Let $A_{S, 1}, \ldots, A_{S, N}$ be the connected components of $A_{S}$. Then $\pi: A_{S, j} \longrightarrow \Delta^{\prime} \backslash S$ is a covering with $q_{j}$ sheets, $q_{1}+\cdots+q_{N}=q$. For every point $\zeta^{\prime} \in \Delta^{\prime} \backslash S$, there exists a neighborhood $U$ of $\zeta^{\prime}$ such that $A_{S, j} \cap \pi^{-1}(U)$ is a disjoint union of graphs $z^{\prime \prime}=g_{j, k}\left(z^{\prime}\right)$ of analytic functions $g_{j, k} \in \mathbb{G}(U), 1 \leqslant k \leqslant q_{j}$. If $\lambda\left(z^{\prime \prime}\right)$ is an arbitrary linear form in $z_{d+1}, \ldots, z_{n}$ and $z^{\prime} \in \Delta^{\prime} \backslash S$, we set

$$
P_{\lambda, j}\left(z^{\prime} ; T\right)=\prod_{\left\{z^{\prime \prime} ;\left(z^{\prime}, z^{\prime \prime}\right) \in A_{S, j}\right\}}\left(T-\lambda\left(z^{\prime \prime}\right)\right)=\prod_{1 \leqslant k \leqslant k_{j}}\left(T-\lambda \circ g_{j, k}\left(z^{\prime}\right)\right) .
$$

This defines a polynomial in $T$ with bounded analytic coefficients on $\Delta^{\prime} \backslash S$. These coefficients have analytic extensions to $\Delta^{\prime}$ (Remark 4.2), thus $P_{\lambda, j} \in \mathscr{O}\left(\Delta^{\prime}\right)[T]$. By construction, $P_{\lambda, j}\left(z^{\prime} ; \lambda\left(z^{\prime \prime}\right)\right)$ vanishes identically on $A_{S, j}$. Set

$$
P_{\lambda}=\prod_{1 \leqslant j \leqslant N} P_{\lambda, j}, \quad f(z)=\delta\left(z^{\prime}\right) P_{\lambda}\left(z^{\prime} ; \lambda\left(z^{\prime \prime}\right)\right)
$$

$f$ vanishes on $A_{S, 1} \cup \ldots \cup A_{S, N} \cup\left(S \times \Delta^{\prime \prime}\right) \supset A \cap \Delta$. Lemma 4.21 shows that $\mathscr{I}_{A, 0}$ is prime; as $\delta \notin \mathscr{I}_{A, 0}$, we get $P_{\lambda, j}\left(z^{\prime} ; \lambda\left(z^{\prime \prime}\right)\right) \in \mathscr{I}_{A, 0}$ for some $j$. This is a contradiction if $N \geqslant 2$ and if $\lambda$ is chosen in such a way that $\lambda$ separates the $q$ points $z_{\nu}^{\prime \prime}$ in each fiber $\pi^{-1}\left(z_{\nu}^{\prime}\right)$, for a sequence $z_{\nu}^{\prime} \rightarrow 0$ in $\Delta^{\prime} \backslash S$. Hence $N=1, A_{S}$ is connected, and for every $\lambda \in\left(\mathbb{C}^{n-d}\right)^{\star}$ we have $P_{\lambda}\left(z^{\prime}, \lambda\left(z^{\prime \prime}\right) \in \mathscr{I}_{(A, 0)}\right.$. By construction $P_{\lambda}\left(z^{\prime}, \lambda\left(z^{\prime \prime}\right)\right)$ vanishes on $A_{S}$, so it vanishes on $\bar{A}_{S}$; hence, for every $z^{\prime} \in S$, the fiber $\bar{A}_{S} \cap \pi^{-1}\left(z^{\prime}\right)$ has at most $q$ elements, otherwise selecting $\lambda$ which separates $q+1$ of these points would yield $q+1$ roots $\lambda\left(z^{\prime \prime}\right)$ of $P_{\lambda}\left(z^{\prime} ; T\right)$, a contradiction. Assume now that $A_{S}$ is not dense in $A \cap \Delta$ for arbitrarily small polydisks $\Delta$. Then there exists a sequence $A \ni z_{\nu}=\left(z_{\nu}^{\prime}, z_{\nu}^{\prime \prime}\right) \rightarrow 0$ such that $z_{\nu}^{\prime} \in S$ and $z_{\nu}^{\prime \prime}$ is not in $F_{\nu}:=\operatorname{pr}^{\prime \prime}\left(\bar{A}_{S} \cap \pi^{-1}\left(z_{\nu}^{\prime}\right)\right)$. The continuity of the roots of the polynomial $P_{\lambda}\left(z^{\prime} ; T\right)$ as $\Delta^{\prime} \backslash S \ni z^{\prime} \rightarrow z_{\nu}^{\prime}$ implies that the set of roots of $P_{\lambda}\left(z_{\nu}^{\prime} ; T\right)$ is $\lambda\left(F_{\nu}\right)$. Select $\lambda$ such that $\lambda\left(z_{\nu}^{\prime \prime}\right) \notin \lambda\left(F_{\nu}\right)$ for all $\nu$. Then $P_{\lambda}\left(z_{\nu}^{\prime} ; \lambda\left(z_{\nu}^{\prime \prime}\right)\right) \neq 0$ for every $\nu$ and $P_{\lambda}\left(z^{\prime} ; \lambda\left(z^{\prime \prime}\right)\right) \notin \mathscr{I}_{A, 0}$, a contradiction.

At this point, it should be observed that many of the above statements completely fail in the case of real analytic sets. In $\mathbb{R}^{2}$, for example, the prime ideal $\mathscr{g}=\left(x^{5}+y^{4}\right)$ defines an irreducible germ of curve $(A, 0)$ and there is an injective integral extension of rings $\mathbb{R}\{x\} \hookrightarrow \mathbb{R}\{x, y\} / \mathscr{G}$ of degree 4 ; however, the projection of $(A, 0)$ on the first factor, $(x, y) \mapsto x$, has not a constant sheet number near 0 , and this number is not related to the degree of the extension. Also, the prime ideal $\mathscr{F}=\left(x^{2}+y^{2}\right)$ has an associated variety $V(\mathscr{F})$ reduced to $\{0\}$, hence $\mathscr{I}_{A, 0}=(x, y)$ is strictly larger than $\mathscr{g}$, in contrast with Lemma 4.21.

Let us return to the complex situation, which is much better behaved. The result obtained in Lemma 4.21 can then be extended to non prime ideals and we get the following important result:
(4.22) Hilbert's Nullstellensatz. For every ideal $\mathscr{G} \subset \mathscr{O}_{n}$

$$
\mathscr{I}_{V(\mathfrak{g}), 0}=\sqrt{\mathfrak{g}},
$$

where $\sqrt{\mathscr{g}}$ is the radical of $\mathscr{g}$, i.e. the set of germs $f \in \mathscr{O}_{n}$ such that some power $f^{k}$ lies in $\mathscr{G}$.

Proof. Set $B=V(\mathscr{y})$. If $f^{k} \in \mathscr{g}$, then $f^{k}$ vanishes on $(B, 0)$ and $f \in \mathscr{I}_{B, 0}$. Thus $\sqrt{\mathscr{g}} \subset \mathscr{I}_{B, 0}$. Conversely, it is well known that $\sqrt{\mathscr{g}}$ is the intersection of all prime ideals $\mathscr{P} \supset \mathscr{J}([$ Lang 1965], Chapter VI). For such an ideal $(B, 0)=(V(\mathscr{y}), 0) \supset(V(\mathscr{P}), 0)$, thus $\mathscr{I}_{B, 0} \subset \mathscr{I}_{V(\mathscr{P}), 0}=\mathscr{P}$ in view of Lemma 4.21. Therefore $\mathscr{I}_{B, 0} \subset \bigcap_{\mathscr{P} \supset \mathscr{g}} \mathscr{P}=\sqrt{\mathscr{J}}$ and the Theorem is proved.

In other words, if a germ $(B, 0)$ is defined by an arbitrary ideal $\mathscr{G} \subset \mathscr{O}_{n}$ and if $f \in \mathscr{O}_{n}$ vanishes on $(B, 0)$, then some power $f^{k}$ lies in $\mathscr{g}$.

## §4.3. Regular and Singular Points. Dimension

The above powerful results enable us to investigate the structure of singularities of an analytic set. We first give a few definitions.
(4.23) Definition. Let $A \subset M$ be an analytic set and $x \in A$. We say that $x \in A$ is a regular point of $A$ if $A \cap \Omega$ is a $\mathbb{C}$-analytic submanifold of $\Omega$ for some neighborhood $\Omega$ of $x$. Otherwise $x$ is said to be singular. The corresponding subsets of $A$ will be denoted respectively $A_{\text {reg }}$ and $A_{\text {sing }}$.

It is clear from the definition that $A_{\text {reg }}$ is an open subset of $A$ (thus $A_{\text {sing }}$ is closed), and that the connected components of $A_{\text {reg }}$ are $\mathbb{C}$-analytic submanifolds of $M$ (non necessarily closed).
(4.24) Proposition. If $(A, x)$ is irreducible, there exist arbitrarily small neighborhoods $\Omega$ of $x$ such that $A_{\mathrm{reg}} \cap \Omega$ is dense and connected in $A \cap \Omega$.

Proof. Take $\Omega=\Delta$ as in Th. 4.19. Then $A_{S} \subset A_{\text {reg }} \cap \Omega \subset A \cap \Omega$, where $A_{S}$ is connected and dense in $A \cap \Omega$; hence $A_{\text {reg }} \cap \Omega$ has the same properties.
(4.25) Definition. The dimension of an irreducible germ of analytic set $(A, x)$ is defined by $\operatorname{dim}(A, x)=\operatorname{dim}\left(A_{\text {reg }}, x\right)$. If $(A, x)$ has several irreducible components $\left(A_{l}, x\right)$, we set

$$
\operatorname{dim}(A, x)=\max \left\{\operatorname{dim}\left(A_{l}, x\right)\right\}, \quad \operatorname{codim}(A, x)=n-\operatorname{dim}(A, x) .
$$

(4.26) Proposition. Let $(B, x) \subset(A, x)$ be germs of analytic sets. If $(A, x)$ is irreducible and $(B, x) \neq(A, x)$, then $\operatorname{dim}(B, x)<\operatorname{dim}(A, x)$ and $B \cap \Omega$ has empty interior in $A \cap \Omega$ for all sufficiently small neighborhoods $\Omega$ of $x$.

Proof. We may assume $x=0,(A, 0) \subset\left(\mathbb{C}^{n}, 0\right)$ and $(B, 0)$ irreducible. Then $\mathscr{I}_{A, 0} \subset \mathscr{I}_{B, 0}$ are prime ideals. When we choose suitable coordinates for the ramified coverings, we may at each step select vectors $e_{n}, e_{n-1}, \ldots$ that work simultaneously for $A$ and $B$. If $\operatorname{dim} B=\operatorname{dim} A$, the process stops for both at the same time, i.e. we get ramified coverings

$$
\pi: A \cap\left(\Delta^{\prime} \times \Delta^{\prime \prime}\right) \longrightarrow \Delta^{\prime}, \quad \pi: B \cap\left(\Delta^{\prime} \times \Delta^{\prime \prime}\right) \longrightarrow \Delta^{\prime}
$$

with ramification loci $S_{A}, S_{B}$. Then $B \cap\left(\left(\Delta^{\prime} \backslash\left(S_{A} \cup S_{B}\right)\right) \times \Delta^{\prime \prime}\right)$ is an open subset of the manifold $A_{S}=A \cap\left(\left(\Delta^{\prime} \backslash S_{A}\right) \times \Delta^{\prime \prime}\right)$, therefore $B \cap A_{S}$ is an analytic subset of $A_{S}$ with non empty interior. The same conclusion would hold if $B \cap \Delta$ had non empty interior in $A \cap \Delta$. As $A_{S}$ is connected, we get $B \cap A_{S}=A_{S}$, and as $B \cap \Delta$ is closed in $\Delta$ we infer $B \cap \Delta \supset \bar{A}_{S}=A \cap \Delta$, hence $(B, 0)=(A, 0)$, in contradiction with the hypothesis.
(4.27) Example: parametrization of curves. Suppose that $(A, 0)$ is an irreducible germ of curve $(\operatorname{dim}(A, 0)=1)$. If the disk $\Delta^{\prime} \subset \mathbb{C}$ is chosen so small that $S=\{0\}$, then $A_{S}$ is a connected covering of $\Delta^{\prime} \backslash\{0\}$ with $q$ sheets. Hence, there exists a covering isomorphism between $\pi$ and the standard covering

$$
\mathbb{C} \supset \Delta(r) \backslash\{0\} \longrightarrow \Delta\left(r^{q}\right) \backslash\{0\}, \quad t \longmapsto t^{q}, \quad r^{q}=\text { radius of } \Delta^{\prime},
$$

i.e. a map $\gamma: \Delta(r) \backslash\{0\} \longrightarrow A_{S}$ such that $\pi \circ \gamma(t)=t^{q}$. This map extends into a bijective holomorphic map $\gamma: \Delta(r) \longrightarrow A \cap \Delta$ with $\gamma(0)=0$. This means that every irreducible germ of curve can be parametrized by a bijective holomorphic map defined on a disk in $\mathbb{C}$ (see also Exercise 10.8).

## $\S$ 4.4. Coherence of Ideal Sheaves

Let $A$ be an analytic set in a complex manifold $M$. The sheaf of ideals $\mathscr{I}_{A}$ is the subsheaf of $\mathscr{O}_{M}$ consisting of germs of holomorphic functions on $M$ which vanish on $A$. Its stalks are the ideals $\mathscr{I}_{A, x}$ already considered; note that $\mathscr{I}_{A, x}=\mathscr{O}_{M, x}$ if $x \notin A$. If $x \in A$, we let $\mathbb{G}_{A, x}$ be the ring of germs of functions on $(A, x)$ which can be extended as germs of holomorphic functions on $(M, x)$. By definition, there is a surjective morphism $\mathfrak{O}_{M, x} \longrightarrow \mathfrak{O}_{A, x}$ whose kernel is $\mathscr{I}_{A, x}$, thus

$$
\begin{equation*}
\mathfrak{O}_{A, x}=\mathfrak{O}_{M, x} / \mathscr{I}_{A, x}, \quad \forall x \in A, \tag{4.28}
\end{equation*}
$$

i.e. $\mathscr{O}_{A}=\left(\mathscr{O}_{M} / \mathscr{I}_{A}\right)_{\mid A}$. Since $\mathscr{I}_{A, x}=\mathscr{O}_{M, x}$ for $x \notin A$, the quotient sheaf $\mathscr{O}_{M} / \mathscr{I}_{A}$ is zero on $M \backslash A$.
(4.29) Theorem ([Cartan 1950]). For any analytic set $A \subset M$, the sheaf of ideals $\mathscr{I}_{A}$ is a coherent analytic sheaf.

Proof. It is sufficient to prove the result when $A$ is an analytic subset in a neighborhood of 0 in $\mathbb{C}^{n}$. If $(A, 0)$ is not irreducible, there exists a neighborhood $\Omega$ such that $A \cap$ $\Omega=A_{1} \cup \ldots \cup A_{N}$ where $A_{k}$ are analytic sets in $\Omega$ and $\left(A_{k}, 0\right)$ is irreducible. We have $\mathscr{I}_{A \cap \Omega}=\bigcap \mathscr{I}_{A_{k}}$, so by Cor. 3.15 we may assume that $(A, 0)$ is irreducible. Then we can choose coordinates $z^{\prime}, z^{\prime \prime}$, polydisks $\Delta^{\prime}, \Delta^{\prime \prime}$ and a primitive element $u\left(z^{\prime \prime}\right)=$ $c_{d+1} z_{d+1}+\cdots+c_{n} z_{n}$ such that Th. 4.19 is valid. Since $\delta\left(z^{\prime}\right)=\prod_{j<k}\left(\sigma_{k} \tilde{u}-\sigma_{j} \tilde{u}\right)^{2}$, we see that $\delta\left(z^{\prime}\right)$ is in fact a polynomial in the $c_{j}$ 's with coefficients in $\mathscr{O}_{d}$. The same is true for the coefficients of the polynomials $W_{u}\left(z^{\prime} ; T\right)$ and $B_{k}\left(z^{\prime} ; T\right)$ which can be expressed in terms of the elementary symmetric functions of the $\sigma_{k} \tilde{u}$ 's. We suppose that $\Delta^{\prime}$ is chosen small enough in order that all coefficients of these $\mathscr{O}_{d}\left[c_{d+1}, \ldots, c_{n}\right]$ polynomials are in $\mathscr{O}\left(\Delta^{\prime}\right)$. Let $\delta_{\alpha} \in \mathscr{O}\left(\Delta^{\prime}\right)$ be some non zero coefficient appearing in $\delta^{m}=\sum \delta_{\alpha} c^{\alpha}$. Also, let $G_{1}, \ldots, G_{N} \in \mathscr{O}\left(\Delta^{\prime}\right)\left[z^{\prime \prime}\right]$ be the coefficients of all monomials $c^{\alpha}$ appearing in the expansion of the functions $W_{u}\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right)$ or $\delta\left(z^{\prime}\right) z_{k}-B_{k}\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right)$. Clearly, $G_{1}, \ldots, G_{N}$ vanish on $A \cap \Delta$. We contend that

$$
\begin{equation*}
\mathscr{I}_{A, x}=\left\{f \in \mathscr{O}_{M, x} ; \delta_{\alpha} f \in\left(G_{1, x}, \ldots, G_{N, x}\right)\right\} . \tag{4.30}
\end{equation*}
$$

This implies that the sheaf $\mathscr{I}_{A}$ is the projection on the first factor of the sheaf of relations $\mathscr{R}\left(\delta_{\alpha}, G_{1}, \ldots, G_{N}\right) \subset \mathscr{O}_{\Delta}^{N+1}$, which is coherent by the Oka theorem; Theorem 4.29 then follows.

We first prove that the inclusion $\mathscr{I}_{A, x} \supset\{\ldots\}$ holds in (4.30). In fact, if $\delta_{\alpha} f \in$ $\left(G_{1, x}, \ldots, G_{N, x}\right)$, then $f$ vanishes on $A \backslash\left\{\delta_{\alpha}=0\right\}$ in some neighborhood of $x$. Since $(A \cap \Delta) \backslash\left\{\delta_{\alpha}=0\right\}$ is dense in $A \cap \Delta$, we conclude that $f \in \mathscr{I}_{A, x}$.

To prove the other inclusion $\mathscr{I}_{A, x} \subset\{\ldots\}$, we repeat the proof of Lemma 4.18 with a few modifications. Let $x \in \Delta$ be a given point. At $x$, the irreducible polynomials $W_{u}\left(z^{\prime} ; T\right)$ and $W_{k}\left(z^{\prime} ; T\right)$ of $\tilde{u}$ and $\tilde{z}_{k}$ in $\mathscr{O}_{M, 0} / \mathscr{I}_{A, 0}$ split into

$$
\begin{aligned}
& W_{u}\left(z^{\prime} ; T\right)=W_{u, x}\left(z^{\prime} ; T-u\left(x^{\prime \prime}\right)\right) Q_{u, x}\left(z^{\prime} ; T-u\left(x^{\prime \prime}\right)\right) \\
& W_{k}\left(z^{\prime} ; T\right)=W_{k, x}\left(z^{\prime} ; T-x_{k}\right) Q_{k, x}\left(z^{\prime} ; T-x_{k}\right)
\end{aligned}
$$

where $W_{u, x}\left(z^{\prime} ; T\right)$ and $W_{k, x}\left(z^{\prime} ; T\right)$ are Weierstrass polynomials in $T$ and $Q_{u, x}\left(x^{\prime}, 0\right) \neq 0$, $Q_{k, x}\left(x^{\prime}, 0\right) \neq 0$. For all $z^{\prime} \in \Delta^{\prime}$, the roots of $W_{u}\left(z^{\prime} ; T\right)$ are the values $u\left(z^{\prime \prime}\right)$ at all points $z \in A \cap \pi^{-1}\left(z^{\prime}\right)$. As $A$ is closed, any point $z \in A \cap \pi^{-1}\left(z^{\prime}\right)$ with $z^{\prime}$ near $x^{\prime}$ has to be in a small neighborhood of one of the points $y \in A \cap \pi^{-1}\left(x^{\prime}\right)$. Choose $c_{d+1}, \ldots, c_{n}$ such that the linear form $u\left(z^{\prime \prime}\right)$ separates all points in the fiber $A \cap \pi^{-1}\left(x^{\prime}\right)$. Then, for a root $u\left(z^{\prime \prime}\right)$ of $W_{u, x}\left(z^{\prime} ; T-u\left(x^{\prime \prime}\right)\right)$, the point $z$ must be in a neighborhood of $y=x$, otherwise $u\left(z^{\prime \prime}\right)$ would be near $u\left(y^{\prime \prime}\right) \neq u\left(x^{\prime \prime}\right)$ and the Weierstrass polynomial $W_{u, x}\left(z^{\prime} ; T\right)$ would have a root away from 0 , in contradiction with (4.10). Conversely, if $z \in A \cap \pi^{-1}\left(z^{\prime}\right)$ is near $x$, then $Q_{u, x}\left(z^{\prime} ; u\left(z^{\prime \prime}\right)-u\left(x^{\prime \prime}\right)\right) \neq 0$ and $u\left(z^{\prime \prime}\right)$ is a root of $W_{u, x}\left(z^{\prime} ; T-u\left(x^{\prime \prime}\right)\right)$. From this, we infer that every polynomial $P\left(z^{\prime} ; T\right) \in \mathbb{O}_{\Delta^{\prime}, x^{\prime}}[T]$ such that $P\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right)=0$ on $(A, x)$ is a multiple of $W_{u, x}\left(z^{\prime} ; T-u\left(x^{\prime \prime}\right)\right)$, because the roots of the latter polynomial are simple for $z^{\prime}$ in the dense set $\left(\Delta^{\prime} \backslash S, x\right)$. In particular $\operatorname{deg} P<\operatorname{deg} W_{u, x}$ implies $P=0$ and

$$
\delta\left(z^{\prime}\right)^{q} W_{k, x}\left(z^{\prime} ; B_{k}\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right) / \delta\left(z^{\prime}\right)-x_{k}\right)
$$

is a multiple of $W_{u, x}\left(z^{\prime} ; T-u\left(x^{\prime \prime}\right)\right)$. If we replace $W_{u}, W_{k}$ by $W_{u, x}, W_{k, x}$ respectively, the proof of Lemma 4.18 shows that for every $f \in \mathscr{O}_{M, x}$ there is a polynomial $R \in \mathscr{O}_{\Delta^{\prime}, x^{\prime}}[T]$ of degree $\operatorname{deg} R<\operatorname{deg} W_{u, x}$ such that

$$
\begin{aligned}
& \delta\left(z^{\prime}\right)^{m} f(z)=R\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right) \quad \text { modulo the ideal } \\
& \left(W_{u, x}\left(z^{\prime} ; u\left(z^{\prime \prime}\right)-u\left(x^{\prime \prime}\right)\right), \delta\left(z^{\prime}\right) z_{k}-B_{k}\left(z^{\prime} ; u\left(z^{\prime \prime}\right)\right)\right),
\end{aligned}
$$

and $f \in \mathscr{I}_{A, x}$ implies $R=0$. Since $W_{u, x}$ differs from $W_{u}$ only by an invertible element in $\mathscr{0}_{M, x}$, we conclude that

$$
\left(\sum \delta_{\alpha} c^{\alpha}\right) \mathscr{I}_{A, x}=\delta^{m} \mathscr{I}_{A, x} \subset\left(G_{1, x}, \ldots, G_{N, x}\right)
$$

This is true for a dense open set of coefficients $c_{d+1}, \ldots, c_{n}$, therefore by expressing the coefficients $\delta_{\alpha}$ through interpolation of $\sum \delta_{\alpha} c^{\alpha}$ at suitable points $c$ we infer

$$
\delta_{\alpha} \mathscr{I}_{A, x} \subset\left(G_{1, x}, \ldots, G_{N, x}\right) \quad \text { for all } \alpha .
$$

(4.31) Theorem. $A_{\text {sing }}$ is an analytic subset of $A$.

Proof. The statement is local. Assume first that $(A, 0)$ is an irreducible germ in $\mathbb{C}^{n}$. Let $g_{1}, \ldots, g_{N}$ be generators of the sheaf $\mathscr{I}_{A}$ on a neighborhood $\Omega$ of 0 . Set $d=\operatorname{dim} A$. In a neighborhood of every point $x \in A_{\text {reg }} \cap \Omega, A$ can be defined by holomorphic equations $u_{1}(z)=\ldots=u_{n-d}(z)=0$ such that $d u_{1}, \ldots, d u_{n-d}$ are linearly independant. As $u_{1}, \ldots, u_{n-d}$ are generated by $g_{1}, \ldots, g_{N}$, one can extract a subfamily $g_{j_{1}}, \ldots, g_{j_{n-d}}$ that has at least one non zero Jacobian determinant of rank $n-d$ at $x$. Therefore $A_{\text {sing }} \cap \Omega$ is defined by the equations

$$
\operatorname{det}\left(\frac{\partial g_{j}}{\partial z_{k}}\right)_{\substack{j \in J \\ k \in K}}=0, \quad J \subset\{1, \ldots, N\}, \quad K \subset\{1, \ldots, n\}, \quad|J|=|K|=n-d
$$

Assume now that $(A, 0)=\bigcup\left(A_{l}, 0\right)$ with $\left(A_{l}, 0\right)$ irreducible. The germ of an analytic set at a regular point is irreducible, thus every point which belongs simultaneously to at least two components is singular. Hence

$$
\left(A_{\text {sing }}, 0\right)=\bigcup\left(A_{l, \text { sing }}, 0\right) \cup \bigcup_{k \neq l}\left(A_{k} \cap A_{l}, 0\right)
$$

and $A_{\text {sing }}$ is analytic.
Now, we give a characterization of regular points in terms of a simple algebraic property of the ring $\mathbb{O}_{A, x}$.
(4.32) Proposition. Let $(A, x)$ be a germ of analytic set of dimensiond and let $\mathfrak{m}_{A, x} \subset$ $\mathcal{O}_{A, x}$ be the maximal ideal of functions that vanish at $x$. Then $\mathfrak{m}_{A, x}$ cannot have less than d generators and $\mathfrak{m}_{A, x}$ has d generators if and only if $x$ is a regular point.

Proof. If $A \subset \mathbb{C}^{n}$ is a $d$-dimensional submanifold in a neighborhood of $x$, there are local coordinates centered at $x$ such that $A$ is given by the equations $z_{d+1}=\ldots=z_{n}$ near $z=0$. Then $\mathbb{O}_{A, x} \simeq \mathfrak{O}_{d}$ and $\mathfrak{m}_{A, x}$ is generated by $z_{1}, \ldots, z_{d}$. Conversely, assume that $\mathfrak{m}_{A, x}$ has $s$ generators $g_{1}(z), \ldots, g_{s}(z)$ in $\mathscr{O}_{A, x}=\mathscr{O}_{\mathbb{C}^{n}, x} / \mathscr{I}_{A, x}$. Letting $x=0$ for simplicity, we can write

$$
z_{j}=\sum_{1 \leqslant k \leqslant s} u_{j k}(z) g_{k}(z)+f_{j}(z), \quad u_{j k} \in \mathscr{O}_{n}, \quad f_{j} \in \mathscr{I}_{A, 0}, \quad 1 \leqslant j \leqslant n .
$$

Then we find $d z_{j}=\sum c_{j k}(0) d g_{k}(0)+d f_{j}(0)$, so that the rank of the system of differentials $\left(d f_{j}(0)\right)_{1 \leqslant j \leqslant n}$ is at least equal to $n-s$. Assume for example that $d f_{1}(0), \ldots, d f_{n-s}(0)$ are linearly independent. By the implicit function theorem, the equations $f_{1}(z)=\ldots=$ $f_{n-s}(z)=0$ define a germ of submanifold of dimension $s$ containing $(A, 0)$, thus $s \geqslant d$ and $(A, 0)$ equals this submanifold if $s=d$.
(4.33) Corollary. Let $A \subset M$ be an analytic set of pure dimension $d$ and let $B \subset A$ be an analytic subset of codimension $\geqslant p$ in $A$. Then, as an $\mathbb{O}_{A, x}$-module, the ideal $\mathscr{I}_{B, x}$ cannot be generated by less than $p$ generators at any point $x \in B$, and by less than $p+1$ generators at any point $x \in B_{\mathrm{reg}} \cap A_{\text {sing }}$.

Proof. Suppose that $\mathscr{I}_{B, x}$ admits $s$-generators $\left(g_{1}, \ldots, g_{s}\right)$ at $x$. By coherence of $\mathscr{I}_{B}$ these germs also generate $\mathscr{I}_{B}$ in a neighborhood of $x$, so we may assume that $x$ is a regular point of $B$. Then there are local coordinates $\left(z_{1}, \ldots, z_{n}\right)$ on $M$ centered at $x$ such that $(B, x)$ is defined by $z_{k+1}=\ldots=z_{n}=0$, where $k=\operatorname{dim}(B, x)$. Then the maximal ideal $\mathfrak{m}_{B, x}=\mathfrak{m}_{A, x} / \mathscr{I}_{B, x}$ is generated by $z_{1}, \ldots, z_{k}$, so that $\mathfrak{m}_{A, x}$ is generated by $\left(z_{1}, \ldots, z_{k}, g_{1}, \ldots, g_{s}\right)$. By Prop. 4.32, we get $k+s \geqslant d$, thus $s \geqslant d-k \geqslant p$, and we have strict inequalities when $x \in A_{\text {sing }}$.

## § 5. Complex Spaces

Much in the same way a manifold is constructed by piecing together open patches isomorphic to open sets in a vector space, a complex space is obtained by gluing together open patches isomorphic to analytic subsets. The general concept of analytic morphism (or holomorphic map between analytic sets) is first needed.

## $\S$ 5.1. Morphisms and Comorphisms

Let $A \subset \Omega \subset \mathbb{C}^{n}$ and $B \subset \Omega^{\prime} \subset \mathbb{C}^{p}$ be analytic sets. A morphism from $A$ to $B$ is by definition a map $F: A \longrightarrow B$ such that for every $x \in A$ there is a neighborhood $U$ of $x$ and a holomorphic map $\tilde{F}: U \longrightarrow \mathbb{C}^{p}$ such that $\tilde{F}_{\mid A \cap U}=F_{\upharpoonright A \cap U}$. Equivalently, such a
morphism can be defined as a continuous map $F: A \longrightarrow B$ such that for all $x \in A$ and $g \in \mathfrak{O}_{B, F(x)}$ we have $g \circ F \in \mathbb{O}_{A, x}$. The induced ring morphism

$$
\begin{equation*}
F_{x}^{\star}: \mathfrak{O}_{B, F(x)} \ni g \longmapsto g \circ F \in \mathbb{O}_{A, x} \tag{5.1}
\end{equation*}
$$

is called the comorphism of $F$ at point $x$.

## $\S$ 5.1. Definition of Complex Spaces

(5.2) Definition. A complex space $X$ is a locally compact Hausdorff space, countable at infinity, together with a sheaf $\mathbb{O}_{X}$ of continuous functions on $X$, such that there exists an open covering $\left(U_{\lambda}\right)$ of $X$ and for each $\lambda$ a homeomorphism $F_{\lambda}: U_{\lambda} \longrightarrow A_{\lambda}$ onto an analytic set $A_{\lambda} \subset \Omega_{\lambda} \subset \mathbb{C}^{n_{\lambda}}$ such that the comorphism $F_{\lambda}^{\star}: \mathscr{O}_{A_{\lambda}} \longrightarrow \mathbb{O}_{X} \mid U_{\lambda}$ is an isomorphism of sheaves of rings. $\mathbb{O}_{X}$ is called the structure sheaf of $X$.

By definition a complex space $X$ is locally isomorphic to an analytic set, so the concepts of holomorphic function on $X$, of analytic subset, of analytic morphism, etc ... are meaningful. If $X$ is a complex space, Th. 4.31 implies that $X_{\text {sing }}$ is an analytic subset of $X$.
(5.3) Theorem and definition. For every complex space $X$, the set $X_{\text {reg }}$ is a dense open subset of $X$, and consists of a disjoint union of connected complex manifolds $X_{\alpha}^{\prime}$. Let $X_{\alpha}$ be the closure of $X_{\alpha}^{\prime}$ in $X$. Then $\left(X_{\alpha}\right)$ is a locally finite family of analytic subsets of $X$, and $X=\bigcup X_{\alpha}$. The sets $X_{\alpha}$ are called the global irreducible components of $X$.


Fig. II-2 The irreducible curve $y^{2}=x^{2}(1+x)$ in $\mathbb{C}^{2}$.

Observe that the germ at a given point of a global irreducible component can be reducible, as shows the example of the cubic curve $\Gamma: y^{2}=x^{2}(1+x)$; the germ $(\Gamma, 0)$ has two analytic branches $y= \pm x \sqrt{1+x}$, however $\Gamma \backslash\{0\}$ is easily seen to be a connected smooth Riemann surface (the real points of $\gamma$ corresponding to $-1 \leqslant x \leqslant 0$ form a path connecting the two branches). This example shows that the notion of global irreducible component is quite different from the notion of local irreducible component introduced in (4.4).

Proof. By definition of $X_{\text {reg }}$, the connected components $X_{\alpha}^{\prime}$ are (disjoint) complex manifolds. Let us show that the germ of $X_{\alpha}=\bar{X}_{\alpha}^{\prime}$ at any point $x \in X$ is analytic. We may
assume that $(X, x)$ is a germ of analytic set $A$ in an open subset of $\mathbb{C}^{n}$. Let $\left(A_{l}, x\right)$, $1 \leqslant l \leqslant N$, be the irreducible components of this germ and $U$ a neighborhood of $x$ such that $X \cap U=\bigcup A_{l} \cap U$. Let $\Omega_{l} \subset U$ be a neighborhood of $x$ such that $A_{l, \text { reg }} \cap \Omega_{l}$ is connected and dense in $A_{l} \cap \Omega_{l}$ (Prop. 4.24). Then $A_{l}^{\prime}:=X_{\text {reg }} \cap A_{l} \cap \Omega_{l}$ equals $\left(A_{l, \text { reg }} \cap \Omega_{l}\right) \backslash \bigcup_{k \neq l} A_{l, \text { reg }} \cap \Omega_{l} \cap A_{k}$. However, $A_{l, \text { reg }} \cap \Omega_{l} \cap A_{k}$ is an analytic subset of $A_{l, \text { reg }} \cap \Omega_{l}$, distinct from $A_{l, \text { reg }} \cap \Omega_{l}$, otherwise $A_{l, \text { reg }} \cap \Omega_{l}$ would be contained in $A_{k}$, thus $\left(A_{l}, x\right) \subset\left(A_{k}, x\right)$ by density. Remark 4.2 implies that $A_{l}^{\prime}$ is connected and dense in $A_{l, \text { reg }} \cap \Omega_{l}$, hence in $A_{l} \cap \Omega_{l}$. Set $\Omega=\bigcap \Omega_{l}$ and let $\left(X_{\alpha}\right)_{\alpha \in J}$ be the family of global components which meet $\Omega$ (i.e. such that $X_{\alpha}^{\prime} \cap \Omega \neq \emptyset$ ). As $X_{\text {reg }} \cap \Omega=\bigcup A_{l}^{\prime} \cap \Omega$, each $X_{\alpha}^{\prime}, \alpha \in J$, meets at least one set $A_{l}^{\prime}$, and as $A_{l}^{\prime} \subset X_{\text {reg }}$ is connected, we have in fact $A_{l}^{\prime} \subset X_{\alpha}^{\prime}$. It follows that there exists a partition $\left(L_{\alpha}\right)_{\alpha \in J}$ of $\{1, \ldots, N\}$ such that $X_{\alpha}^{\prime} \cap \Omega=\bigcup_{l \in L_{\alpha}} A_{l}^{\prime} \cap \Omega, \alpha \in J$. Hence $J$ is finite, card $J \leqslant N$, and

$$
X_{\alpha} \cap \Omega=\bar{X}_{\alpha}^{\prime} \cap \Omega=\bigcup_{l \in L_{\alpha}} \bar{A}_{l}^{\prime} \cap \Omega=\bigcup_{l \in L_{\alpha}} A_{l} \cap \Omega
$$

is analytic for all $\alpha \in J$.
(5.4) Corollary. If $A, B$ are analytic subsets in a complex space $X$, then the closure $\overline{A \backslash B}$ is an analytic subset, consisting of the union of all global irreducible components $A_{\lambda}$ of $A$ which are not contained in $B$.

Proof. Let $C=\bigcup A_{\lambda}$ be the union of these components. Since $\left(A_{\lambda}\right)$ is locally finite, $C$ is analytic. Clearly $A \backslash B=C \backslash B=\bigcup A_{\lambda} \backslash B$. The regular part $A_{\lambda}^{\prime}$ of each $A_{\lambda}$ is a connected manifold and $A_{\lambda}^{\prime} \cap B$ is a proper analytic subset (otherwise $A_{\lambda}^{\prime} \subset B$ would imply $\left.A_{\lambda} \subset B\right)$. Thus $A_{\lambda}^{\prime} \backslash\left(A_{\lambda}^{\prime} \cap B\right)$ is dense in $A_{\lambda}^{\prime}$ which is dense in $A_{\lambda}$, so $\overline{A \backslash B}=\bigcup A_{\lambda}=C$.
(5.5) Theorem. For any family $\left(A_{\lambda}\right)$ of analytic sets in a complex space $X$, the intersection $A=\bigcap A_{\lambda}$ is an analytic subset of $X$. Moreover, the intersection is stationary on every compact subset of $X$.

Proof. It is sufficient to prove the last statement, namely that every point $x \in X$ has a neighborhood $\Omega$ such that $A \cap \Omega$ is already obtained as a finite intersection. However, since $\mathbb{O}_{X, x}$ is Noetherian, the family of germs of finite intersections has a minimum element $(B, x), B=\bigcap A_{\lambda_{j}}, 1 \leqslant j \leqslant N$. Let $\tilde{B}$ be the union of the global irreducible components $B_{\alpha}$ of $B$ which contain the point $x$; clearly $(B, x)=(\tilde{B}, x)$. For any set $A_{\lambda}$ in the family, the minimality of $B$ implies $(B, x) \subset\left(A_{\lambda}, x\right)$. Let $B_{\alpha}^{\prime}$ be the regular part of any global irreducible component $B_{\alpha}$ of $\tilde{B}$. Then $B_{\alpha}^{\prime} \cap A_{\lambda}$ is a closed analytic subset of $B_{\alpha}^{\prime}$ containing a non empty open subset (the intersection of $B_{\alpha}^{\prime}$ with some neighborhood of $x$ ), so we must have $B_{\alpha}^{\prime} \cap A_{\lambda}=B_{\alpha}^{\prime}$. Hence $B_{\alpha}=\bar{B}_{\alpha}^{\prime} \subset A_{\lambda}$ for all $B_{\alpha} \subset \tilde{B}$ and all $A_{\lambda}$, thus $\tilde{B} \subset A=\bigcap A_{\lambda}$. We infer

$$
(B, x)=(\tilde{B}, x) \subset(A, x) \subset(B, x)
$$

and the proof is complete.
As a consequence of these general results, it is not difficult to show that a complex space always admits a (locally finite) stratification such that the strata are smooth manifolds.
(5.6) Proposition. Let $X$ be a complex space. Then there is a locally stationary increasing sequence of analytic subsets $Y_{k} \subset X, k \in \mathbb{N}$, such that $Y_{0}$ is a discrete set and such that $Y_{k} \backslash Y_{k-1}$ is a smooth $k$-dimensional complex manifold for $k \geqslant 1$. Such a sequence is called a stratification of $X$, and the sets $Y_{k} \backslash Y_{k-1}$ are called the strata (the strata may of course be empty for some indices $k<\operatorname{dim} X$ ).

Proof. Let $\mathscr{F}$ be the family of irreducible analytic subsets $Z \subset X$ which can be obtained through a finite sequence of steps of the following types:
a) $Z$ is an irreducible component of $X$;
b) $Z$ is an irreducible component of the singular set $Z_{\text {sing }}^{\prime}$ of some member $Z^{\prime} \in \mathscr{F}$;
c) $Z$ is an irreducible component of some finite intersection of sets $Z_{j} \in \mathscr{F}$.

Since $X$ has locally finite dimension and since steps b) or c) decrease the dimension of our sets $Z$, it is clear that $\mathscr{F}$ is a locally finite family of analytic sets in $X$. Let $Y_{k}$ be the union of all sets $Z \in \mathscr{F}$ of dimension $\leqslant k$. It is easily seen that $\bigcup Y_{k}=X$ and that the irreducible components of $\left(Y_{k}\right)_{\text {sing }}$ are contained in $Y_{k-1}$ (these components are either intersections of components $Z_{j} \subset Y_{k}$ or parts of the singular set of some component $Z \subset Y_{k}$, so there are in either case obtained by step b) or c) above). Hence $Y_{k} \backslash Y_{k-1}$ is a smooth manifold and it is of course $k$-dimensional, because the components of $Y_{k}$ of dimension $<k$ are also contained in $Y_{k-1}$ by definition.
(5.7) Theorem. Let $X$ be an irreducible complex space. Then every non constant holomorphic function $f$ on $X$ defines an open map $f: X \longrightarrow \mathbb{C}$.

Proof. We show that the image $f(\Omega)$ of any neighborhood $\Omega$ of $x \in X$ contains a neighborhood of $f(x)$. Let $\left(X_{l}, x\right)$ be an irreducible component of the germ $(X, x)$ (embedded in $\mathbb{C}^{n}$ ) and $\Delta=\Delta^{\prime} \times \Delta^{\prime \prime} \subset \Omega$ a polydisk such that the projection $\pi: X_{l} \cap \Delta \longrightarrow \Delta^{\prime}$ is a ramified covering. The function $f$ is non constant on the dense open manifold $X_{\text {reg }}$, so we may select a complex line $L \subset \Delta^{\prime}$ through 0 , not contained in the ramification locus of $\pi$, such that $f$ is non constant on the one dimensional germ $\pi^{-1}(L)$. Therefore we can find a germ of curve

$$
(\mathbb{C}, 0) \ni t \longmapsto \gamma(t) \in(X, x)
$$

such that $f \circ \gamma$ is non constant. This implies that the image of every neighborhood of $0 \in \mathbb{C}$ by $f \circ \gamma$ already contains a neighborhood of $f(x)$.
(5.8) Corollary. If $X$ is a compact irreducible analytic space, then every holomorphic function $f \in \mathscr{O}(X)$ is constant.

In fact, if $f \in \mathscr{O}(X)$ was non constant, $f(X)$ would be compact and also open in $\mathbb{C}$ by Th. 5.7, a contradiction. This result implies immediately the following consequence.
(5.9) Theorem. Let $X$ be a complex space such that the global holomorphic functions in $\mathscr{G}(X)$ separate the points of $X$. Then every compact analytic subset $A$ of $X$ is finite.

Proof. $A$ has a finite number of irreducible components $A_{\lambda}$ which are also compact. Every function $f \in \mathscr{O}(X)$ is constant on $A_{\lambda}$, so $A_{\lambda}$ must be reduced to a single point.

## $\S$ 5.2. Coherent Sheaves over Complex Spaces

Let $X$ be a complex space and $\mathscr{O}_{X}$ its structure sheaf. Locally, $X$ can be identified to an analytic set $A$ in an open set $\Omega \subset \mathbb{C}^{n}$, and we have $\mathscr{O}_{X}=\mathscr{O}_{\Omega} / \mathscr{F}_{A}$. Thus $\mathscr{O}_{X}$ is coherent over the sheaf of rings $\mathscr{O}_{\Omega}$. It follows immediately that $\mathbb{O}_{X}$ is coherent over itself. Let $\mathscr{S}$ be a $\mathbb{O}_{X}$-module. If $\mathscr{S}$ denotes the extension of $\mathscr{S}_{\upharpoonright A}$ to $\Omega$ obtained by setting $\tilde{\mathscr{S}}_{x}=0$ for $x \in \Omega \backslash A$, then $\tilde{\mathscr{T}}$ is a $\mathscr{O}_{\Omega}$-module, and it is easily seen that $\mathscr{\mathscr { S }}_{\uparrow A}$ is coherent over $\mathscr{O}_{X \upharpoonright A}$ if and only if $\tilde{\mathscr{S}}$ is coherent over $\mathscr{O}_{\Omega}$. If $Y$ is an analytic subset of $X$, then $Y$ is locally given by an analytic subset $B$ of $A$ and the sheaf of ideals of $Y$ in $\mathbb{O}_{X}$ is the quotient $\mathscr{I}_{Y}=\mathscr{I}_{B} / \mathscr{I}_{A}$; hence $\mathscr{I}_{Y}$ is coherent. Let us mention the following important property of supports.
(5.10) Theorem. If $\mathscr{S}$ is a coherent $\mathscr{G}_{X}$-module, the support of $\mathscr{S}$, defined as $\operatorname{Supp} \mathscr{\mathscr { S }}=$ $\left\{x \in X ; \mathscr{S}_{x} \neq 0\right\}$ is an analytic subset of $X$.

Proof. The result is local, thus after extending $\mathscr{S}$ by 0 , we may as well assume that $X$ is an open subset $\Omega \subset \mathbb{C}^{n}$. By (3.12), there is an exact sequence of sheaves

$$
\mathfrak{O}_{U}^{\oplus p} \xrightarrow{G} \mathscr{O}_{U}^{\oplus q} \xrightarrow{F} \mathscr{S}_{\upharpoonright U} \longrightarrow 0
$$

in a neighborhood $U$ of any point. If $G: \mathscr{O}_{x}^{\oplus p} \longrightarrow \mathscr{O}_{x}^{\oplus q}$ is surjective it is clear that the linear map $G(x): \mathbb{C}^{p} \longrightarrow \mathbb{C}^{q}$ must be surjective; conversely, if $G(x)$ is surjective, there is a $q$-dimensional subspace $E \subset \mathbb{C}^{p}$ on which the restriction of $G(x)$ is a bijection onto $\mathbb{C}^{q}$; then $G_{\upharpoonright E}: \mathscr{O}_{U} \otimes_{\mathbb{C}} E \longrightarrow \mathscr{O}_{U}^{\oplus q}$ is bijective near $x$ and $G$ is surjective. The support of $\mathscr{S}_{\mid U}$ is thus equal to the set of points $x \in U$ such that all minors of $G(x)$ of order $q$ vanish.

## $\S$ 6. Analytic Cycles and Meromorphic Functions

## § 6.1. Complete Intersections

Our goal is to study in more details the dimension of a subspace given by a set of equations. The following proposition is our starting point.
(6.1) Proposition. Let $X$ be a complex space of pure dimension $p$ and $A$ an analytic subset of $X$ with $\operatorname{codim}_{X} A \geqslant 2$. Then every function $f \in \mathbb{O}(X \backslash A)$ is locally bounded near $A$.

Proof. The statement is local on $X$, so we may assume that $X$ is an irreducible germ of analytic set in $\left(\mathbb{C}^{n}, 0\right)$. Let $\left(A_{k}, 0\right)$ be the irreducible components of $(A, 0)$. By a reasoning analogous to that of Prop. 4.26 , we can choose coordinates $\left(z_{1}, \ldots, z_{n}\right)$ on $\mathbb{C}^{n}$ such that all projections

$$
\begin{gathered}
\pi: z \longmapsto\left(z_{1}, \ldots, z_{p}\right), \quad p=\operatorname{dim} X, \\
\pi_{k}: z \longmapsto\left(z_{1}, \ldots, z_{p_{k}}\right), \quad p_{k}=\operatorname{dim} A_{k},
\end{gathered}
$$

define ramified coverings $\pi: X \cap \Delta \longrightarrow \Delta^{\prime}, \pi_{k}: A_{k} \cap \Delta \longrightarrow \Delta_{k}^{\prime}$. By construction $\pi\left(A_{k}\right) \subset \Delta^{\prime}$ is contained in the set $B_{k}$ defined by some Weierstrass polynomials in the
variables $z_{p_{k}+1}, \ldots, z_{p}$ and $\operatorname{codim}_{\Delta^{\prime}} B_{k}=p-p_{k} \geqslant 2$. Let $S$ be the ramification locus of $\pi$ and $B=\bigcup B_{k}$. We have $\pi(A \cap \Delta) \subset B$. For $z^{\prime} \in \Delta^{\prime} \backslash(S \cup B)$, we let

$$
\sigma_{k}\left(z^{\prime}\right)=\text { elementary symmetric function of degree } k \text { in } f\left(z^{\prime}, z_{\alpha}^{\prime \prime}\right),
$$

where $\left(z^{\prime}, z_{\alpha}^{\prime \prime}\right)$ are the $q$ points of $X$ projecting on $z^{\prime}$. Then $\sigma_{k}$ is holomorphic on $\Delta^{\prime} \backslash$ $(S \cup B)$ and locally bounded near every point of $S \backslash B$, thus $\sigma_{k}$ extends holomorphically to $\Delta^{\prime} \backslash B$ by Remark 4.2. Since codim $B \geqslant 2, \sigma_{k}$ extends to $\Delta^{\prime}$ by Cor. 1.4.5. Now, $f$ satisfies $f^{q}-\sigma_{1} f^{q-1}+\ldots+(-1)^{q} \sigma_{q}=0$, thus $f$ is locally bounded on $X \cap \Delta$.
(6.2) Theorem. Let $X$ be an irreducible complex space and $f \in \mathscr{O}(X), f \not \equiv 0$. Then $f^{-1}(0)$ is empty or of pure dimension $\operatorname{dim} X-1$.

Proof. Let $A=f^{-1}(0)$. By Prop. 4.26, we know that $\operatorname{dim} A \leqslant \operatorname{dim} X-1$. If $A$ had an irreducible branch $A_{j}$ of dimension $\leqslant \operatorname{dim} X-2$, then in virtue of Prop. 6.1 the function $1 / f$ would be bounded in a neighborhood of $A_{j} \backslash \bigcup_{k \neq j} A_{k}$, a contradiction.
(6.3) Corollary. If $f_{1}, \ldots, f_{p}$ are holomorphic functions on an irreducible complex space $X$, then all irreducible components of $f_{1}^{-1}(0) \cap \ldots \cap f_{p}^{-1}(0)$ have codimension $\geqslant p$.
(6.4) Definition. Let $X$ be a complex space of pure dimension $n$ and $A$ an analytic subset of $X$ of pure dimension. Then $A$ is said to be a local (set theoretic) complete intersection in $X$ if every point of $A$ has a neighborhood $\Omega$ such that

$$
A \cap \Omega=\left\{x \in \Omega ; f_{1}(x)=\ldots=f_{p}(x)=0\right\}
$$

with exactly $p=\operatorname{codim} A$ functions $f_{j} \in \mathscr{O}(\Omega)$.
(6.5) Remark. As a converse to Th. 6.2, one may ask whether every hypersurface $A$ in $X$ is locally defined by a single equation $f=0$. In general the answer is negative. A simple counterexample for $\operatorname{dim} X=3$ is obtained with the singular quadric $X=$ $\left\{z_{1} z_{2}+z_{3} z_{4}=0\right\} \subset \mathbb{C}^{4}$ and the plane $A=\left\{z_{1}=z_{3}=0\right\} \subset X$. Then $A$ cannot be defined by a single equation $f=0$ near the origin, otherwise the plane $B=\left\{z_{2}=z_{4}=0\right\}$ would be such that

$$
f^{-1}(0) \cap B=A \cap B=\{0\},
$$

in contradiction with Th. 6.2 (also, by Exercise 10.11, we would get the inequality $\operatorname{codim}_{X} A \cap B \leqslant 2$ ). However, the answer is positive when $X$ is a manifold:
(6.6) Theorem. Let $M$ be a complex manifold with $\operatorname{dim}_{\mathbb{C}} M=n$, let $(A, x)$ be an analytic germ of pure dimension $n-1$ and let $A_{j}, 1 \leqslant j \leqslant N$, be its irreducible components.
a) The ideal of $(A, x)$ is a principal ideal $\mathscr{I}_{A, x}=(g)$ where $g$ is a product of irreducible germs $g_{j}$ such that $\mathscr{I}_{A_{j}, x}=\left(g_{j}\right)$.
b) For every $f \in \mathscr{O}_{M, x}$ such that $f^{-1}(0) \subset(A, x)$, there is a unique decomposition $f=u g_{1}^{m_{1}} \ldots g_{N}^{m_{N}}$ where $u$ is an invertible germ and $m_{j}$ is the order of vanishing of $f$ at any point $z \in A_{j, \text { reg }} \backslash \bigcup_{k \neq j} A_{k}$.

Proof. a) In a suitable local coordinate system centered at $x$, the projection $\pi: \mathbb{C}^{n} \longrightarrow$ $\mathbb{C}^{n-1}$ realizes all $A_{j}$ as ramified coverings

$$
\pi: A_{j} \cap \Delta \longrightarrow \Delta^{\prime} \subset \mathbb{C}^{n-1}, \quad \text { ramification locus }=S_{j} \subset \Delta^{\prime}
$$

The function

$$
g_{j}\left(z^{\prime}, z_{n}\right)=\prod_{w \in A_{j} \cap \pi^{-1}\left(z^{\prime}\right)}\left(z_{n}-w_{n}\right), \quad z^{\prime} \in \Delta^{\prime} \backslash S_{j}
$$

extends into a holomorphic function in $\mathbb{O}_{\Delta^{\prime}}\left[z_{n}\right]$ and is irreducible at $x$. Set $g=\prod g_{j} \in$ $\mathscr{I}_{A, x}$. For any $f \in \mathscr{I}_{A, x}$, the Weierstrass division theorem yields $f=g Q+R$ with $R \in \mathbb{O}_{n-1}\left[z_{n}\right]$ and $\operatorname{deg} R<\operatorname{deg} g$. As $R\left(z^{\prime}, z_{n}\right)$ vanishes when $z_{n}$ is equal to $w_{n}$ for each point $w \in A \cap \pi^{-1}\left(z^{\prime}\right), R$ has exactly $\operatorname{deg} g$ roots when $z^{\prime} \in \Delta^{\prime} \backslash\left(\bigcup S_{j} \cup \bigcup \pi\left(A_{j} \cap A_{k}\right)\right)$, so $R=0$. Hence $\mathscr{I}_{A, x}=(g)$ and similarly $\mathscr{I}_{A_{j}, x}=\left(g_{j}\right)$. Since $\mathscr{I}_{A_{j}}$ is coherent, $g_{j}$ is also a generator of $\mathscr{I}_{A_{j}, z}$ for $z$ near $x$ and we infer that $g_{j}$ has order 1 at any regular point $z \in A_{j, \text { reg }}$.
b) As $\mathscr{O}_{M, x}$ is factorial, any $f \in \mathscr{O}_{M, x}$ can be written $f=u g_{1}^{m_{1}} \ldots g_{N}^{m_{N}}$ where $u$ is either invertible or a product of irreducible elements distinct from the $g_{j}$ 's. In the latter case the hypersurface $u^{-1}(0)$ cannot be contained in $(A, x)$, otherwise it would be a union of some of the components $A_{j}$ and $u$ would be divisible by some $g_{j}$. This proves b).
(6.7) Definition. Let $X$ be an complex space of pure dimension $n$.
a) An analytic $q$-cycle $Z$ on $X$ is a formal linear combination $\sum \lambda_{j} A_{j}$ where $\left(A_{j}\right)$ is a locally finite family of irreducible analytic sets of dimension $q$ in $X$ and $\lambda_{j} \in \mathbb{Z}$. The support of $Z$ is $|Z|=\bigcup_{\lambda_{j} \neq 0} A_{j}$. The group of all $q$-cycles on $X$ is denoted $\operatorname{Cycl}^{q}(X)$. Effective $q$-cycles are elements of the subset $\operatorname{Cycl}_{+}^{q}(X)$ of cycles such that all coefficients $\lambda_{j}$ are $\geqslant 0$; rational, real cycles are cycles with coefficients $\lambda_{j} \in \mathbb{Q}, \mathbb{R}$.
b) An analytic $(n-1)$-cycle is called a (Weil) divisor, and we set

$$
\operatorname{Div}(X)=\operatorname{Cycl}^{n-1}(X)
$$

c) Assume that $\operatorname{dim} X_{\text {sing }} \leqslant n-2$. If $f \in \mathscr{O}(X)$ does not vanish identically on any irreducible component of $X$, we associate to $f$ a divisor

$$
\operatorname{div}(f)=\sum m_{j} A_{j} \in \operatorname{Div}_{+}(X)
$$

in the following way: the components $A_{j}$ are the irreducible components of $f^{-1}(0)$ and the coefficient $m_{j}$ is the vanishing order of $f$ at every regular point in $X_{\mathrm{reg}} \cap$ $A_{j, \mathrm{reg}} \backslash \bigcup_{k \neq j} A_{k}$. It is clear that we have

$$
\operatorname{div}(f g)=\operatorname{div}(f)+\operatorname{div}(g)
$$

d) A Cartier divisor is a divisor $D=\sum \lambda_{j} A_{j}$ that is equal locally to a $\mathbb{Z}$-linear combination of divisors of the form $\operatorname{div}(f)$.

It is easy to check that the collection of abelian groups $\operatorname{Cycl}^{q}(U)$ over all open sets $U \subset X$, together with the obvious restriction morphisms, satisfies axioms (1.4) of sheaves; observe however that the restriction of an irreducible component $A_{j}$ to a smaller open set may subdivide in several components. Hence we obtain sheaves of abelian groups $\mathrm{Cycl}^{q}$ and Div $=\mathrm{Cycl}^{n-1}$ on $X$. The stalk $\mathrm{Cycl}_{x}^{q}$ is the free abelian group generated by the set of irreducible germs of $q$-dimensional analytic sets at the point $x$. These sheaves
carry a natural partial ordering determined by the subsheaf of positive elements $\mathrm{Cycl}_{+}^{q}$. We define the sup and inf of two analytic cycles $Z=\sum \lambda_{j} A_{j}, Z^{\prime}=\sum \mu_{j} A_{j}$ by

$$
\begin{equation*}
\sup \left\{Z, Z^{\prime}\right\}=\sum \sup \left\{\lambda_{j}, \mu_{j}\right\} A_{j}, \quad \inf \left\{Z, Z^{\prime}\right\}=\sum \inf \left\{\lambda_{j}, \mu_{j}\right\} A_{j} \tag{6.8}
\end{equation*}
$$

it is clear that these operations are compatible with restrictions, i.e. they are defined as sheaf operations.
(6.9) Remark. When $X$ is a manifold, Th. 6.6 shows that every effective $\mathbb{Z}$-divisor is locally the divisor of a holomorphic function; thus, for manifolds, the concepts of Weil and Cartier divisors coincide. This is not always the case in general: in Example 6.5, one can show that $A$ is not a Cartier divisor (exercise 10.?).

## §6.2. Divisors and Meromorphic Functions

Let $X$ be a complex space. For $x \in X$, let $\mathcal{M}_{X, x}$ be the ring of quotients of $\mathscr{O}_{X, x}$, i.e. the set of formal quotients $g / h, g, h \in \mathbb{O}_{X, x}$, where $h$ is not a zero divisor in $\mathbb{O}_{X, x}$, with the identification $g / h=g^{\prime} / h^{\prime}$ if $g h^{\prime}=g^{\prime} h$. We consider the disjoint union

$$
\begin{equation*}
\mathcal{M}_{X}=\coprod_{x \in X} \mathcal{M}_{X, x} \tag{6.10}
\end{equation*}
$$

with the topology in which the open sets open sets are unions of sets of the type $\left\{G_{x} / H_{x} ; x \in V\right\} \subset \mathcal{M}_{X}$ where $V$ is open in $X$ and $G, H \in \mathbb{O}_{X}(V)$. Then $\mathcal{M}_{X}$ is a sheaf over $X$, and the sections of $\mathscr{M}_{X}$ over an open set $U$ are called meromorphic functions on $U$. By definition, these sections can be represented locally as quotients of holomorphic functions, but there need not exist such a global representation on $U$.

A point $x \in X$ is called a pole of a meromorphic function $f$ on $X$ if $f_{x} \notin \mathscr{O}_{X, x}$. Clearly, the set $P_{f}$ of poles of $f$ is a closed subset of $X$ with empty interior: if $f=g / h$ on $U$, then $h \not \equiv 0$ on any irreducible component and $P_{f} \cap U \subset h^{-1}(0)$. For $x \notin P_{f}$, one can speak of the value $f(x)$. If the restriction of $f$ to $X_{\text {reg }} \backslash P_{f}$ does not vanish identically on any irreducible component of $(X, x)$, then $1 / f$ is a meromorphic function in a neighborhood of $x$; the set of poles of $1 / f$ will be denoted $Z_{f}$ and called the zero set of $f$. If $f$ vanishes on some connected open subset of $X_{\text {reg }} \backslash P_{f}$, then $f$ vanishes identically (outside $P_{f}$ ) on the global irreducible component $X_{\alpha}$ containing this set; we agree that these components $X_{\alpha}$ are contained in $Z_{f}$. For every point $x$ in the complement of $Z_{f} \cap P_{f}$, we have either $f_{x} \in \mathbb{O}_{X, x}$ or $(1 / f)_{x} \in \mathbb{O}_{X, x}$, thus $f$ defines a holomorphic map $X \backslash\left(Z_{f} \cap P_{f}\right) \longrightarrow \mathbb{C} \cup\{\infty\}=\mathbb{P}^{1}$ with values in the projective line. In general, no value (finite or infinite) can be assigned to $f$ at a point $x \in Z_{f} \cap P_{f}$, as shows the example of the function $f(z)=z_{2} / z_{1}$ in $\mathbb{C}^{2}$. The set $Z_{f} \cap P_{f}$ is called the indeterminacy set of $f$.
(6.11) Theorem. For every meromorphic function $f$ on $X$, the sets $P_{f}, Z_{f}$ and the indeterminacy set $Z_{f} \cap P_{f}$ are analytic subsets.

Proof. Let $\mathscr{g}_{x}$ be the ideal of germs $u \in \mathscr{O}_{X, x}$ such that $u f_{x} \in \mathscr{O}_{X, x}$. Let us write $f=g / h$ on a small open set $U$. Then $\mathscr{V}_{V U}$ appears as the projection on the first factor of the sheaf of relations $\mathscr{R}(g, h) \subset \mathfrak{O}_{U} \times \mathscr{O}_{U}$, so $\mathscr{g}$ is a coherent sheaf of ideals. Now

$$
P_{f}=\left\{x \in X ; \mathscr{g}_{x}=\mathscr{O}_{X, x}\right\}=\text { Supp } \mathscr{O}_{X} / \mathscr{f},
$$

thus $P_{f}$ is analytic by Th. 5.10. Similarly, the projection of $\mathscr{R}(g, h)$ on the second factor defines a sheaf of ideals $\mathscr{g}^{\prime}$ such that $Z_{f}=\operatorname{Supp} \mathscr{O}_{X} / \mathscr{g}^{\prime}$.

When $X$ has pure dimension $n$ and $\operatorname{dim} X_{\text {sing }} \leqslant n-2$, Def. 6.7 c) can be extended to meromorphic functions: if $f=g / h$ locally, we set

$$
\begin{equation*}
\operatorname{div}(f)=\operatorname{div}(g)-\operatorname{div}(h) . \tag{6.12}
\end{equation*}
$$

By 6.7 c), we immediately see that this definition does not depend on the choice of the local representant $g / h$. Furthermore, Cartier divisors are precisely those divisors which are associated locally to meromorphic functions.

Assume from now on that $M$ is a connected $n$-dimensional complex manifold. Then, for every point $x \in M$, the ring $\mathbb{O}_{M, x} \simeq \mathscr{O}_{n}$ is factorial. This property makes the study of meromorphic functions much easier.
(6.13) Theorem. Let $f$ be a non zero meromorphic function on a manifold $M$, with $\operatorname{dim}_{\mathbb{C}} M=n$. Then the sets $Z_{f}, P_{f}$ are purely ( $n-1$ )-dimensional, and the indeterminacy set $Z_{f} \cap P_{f}$ is purely ( $n-2$ )-dimensional.

Proof. For every point $a \in M$, the germ $f_{a}$ can be written $g_{a} / h_{a}$ where $g_{a}, h_{a} \in \mathbb{O}_{M, a}$ are relatively prime holomorphic germs. By Th. 1.12, the germs $g_{x}, h_{x}$ are still relatively prime for $x$ in a neighborhood $U$ of $a$. Thus the ideal $\mathscr{f}$ associated to $f$ coincides with ( $h$ ) on $U$, and we have

$$
P_{f} \cap U=\operatorname{Supp} \mathscr{O}_{U} /(h)=h^{-1}(0), \quad Z_{f} \cap U=g^{-1}(0) .
$$

Th. 6.2 implies our contentions: if $g_{\lambda}$ and $h_{\mu}$ are the irreducible components of $g, h$, then $Z_{f} \cap P_{f}=\bigcup g_{\lambda}^{-1}(0) \cap h_{\mu}^{-1}(0)$ is $(n-2)$-dimensional. As we will see in the next section, Th. 6.13 does not hold on an arbitrary complex space.

Let $\left(A_{j}\right)$, resp. $\left(B_{j}\right)$, be the global irreducible components of $Z_{f}$, resp. $P_{f}$. In a neighborhood $V_{j}$ of the $(n-1)$-dimensional analytic set

$$
A_{j}^{\prime}=A_{j} \backslash\left(P_{f} \cup \bigcup_{k \neq j} A_{k}\right)
$$

$f$ is holomorphic and $V \cap f^{-1}(0)=A_{j}^{\prime}$. As $A_{j, \text { reg }}^{\prime}$ is connected, we must have $\operatorname{div}\left(f_{\mid V_{j}}\right)=$ $m_{j} A_{j}^{\prime}$ for some constant multiplicity $m_{j}$ equal to the vanishing order of $f$ along $A_{j, \text { reg }}^{\prime}$. Similarly, $1 / f$ is holomorphic in a neighborhood $W_{j}$ of

$$
B_{j}^{\prime}=B_{j} \backslash\left(Z_{f} \cup \bigcup_{k \neq j} B_{k}\right)
$$

and we have $\operatorname{div}\left(f_{\mid V}\right)=-p_{j} B_{j}^{\prime}$ where $p_{j}$ is the vanishing order of $1 / f$ along $B_{j, \text { reg }}^{\prime}$. At a point $x \in M$ the germs $A_{j, x}$ and $B_{j, x}$ may subdivide in irreducible local components $A_{j, \lambda, x}$ and $B_{j, \lambda, x}$. If $g_{j, \lambda}$ and $h_{j, \lambda}$ are local generators of the corresponding ideals, we may a priori write

$$
f_{x}=u g / h \quad \text { where } \quad g=\prod g_{j, \lambda}^{m_{j, \lambda}}, \quad h=\prod h_{j, \lambda}^{p_{j, \lambda}}
$$

and where $u$ is invertible. Then necessarily $m_{j, \lambda}=m_{j}$ and $p_{j, \lambda}=p_{j}$ for all $\lambda$, and we see that the global divisor of $f$ on $M$ is

$$
\begin{equation*}
\operatorname{div}(f)=\sum m_{j} A_{j}-\sum p_{j} B_{j} . \tag{6.14}
\end{equation*}
$$

Let us denote by $\mathscr{M}^{\star}$ the multiplicative sheaf of germs of non zero meromorphic functions, and by $\mathscr{O}^{\star}$ the sheaf of germs of invertible holomorphic functions. Then we have an exact sequence of sheaves

$$
\begin{equation*}
1 \longrightarrow \mathscr{O}^{\star} \longrightarrow \mathscr{M}^{\star} \xrightarrow{\text { div }} \operatorname{Div} \longrightarrow 0 \tag{6.15}
\end{equation*}
$$

Indeed, the surjectivity of div is a consequence of Th. 6.6. Moreover, any meromorphic function that has a positive divisor must be holomorphic by the fact that $\mathscr{O}_{n}$ is factorial. Hence a meromorphic function $f$ with $\operatorname{div}(f)=0$ is an invertible holomorphic function.

## § 7. Normal Spaces and Normalization

## $\S$ 7.1. Weakly Holomorphic Functions

The goal of this section is to show that the singularities of $X$ can be studied by enlarging the structure sheaf $\mathscr{O}_{X}$ into a sheaf $\tilde{\mathscr{O}}_{X}$ of so-called weakly holomorphic functions.
(7.1) Definition. Let $X$ be a complex space. A weakly holomorphic function $f$ on $X$ is a holomorphic function on $X_{\text {reg }}$ such that every point of $X_{\text {sing }}$ has a neighborhood $V$ for which $f$ is bounded on $X_{\mathrm{reg}} \cap V$. We denote by $\tilde{\mathscr{O}}_{X, x}$ the ring of germs of weakly holomorphic functions over neighborhoods of $x$ and $\tilde{\mathscr{O}}_{X}$ the associated sheaf.

Clearly, $\tilde{\mathscr{O}}_{X, x}$ is a ring containing $\mathscr{C}_{X, x}$. If $\left(X_{j}, x\right)$ are the irreducible components of ( $X, x$ ), there is a fundamental system of neighborhoods $V$ of $x$ such that $X_{\text {reg }} \cap V$ is a disjoint union of connected open sets

$$
X_{j, \mathrm{reg}} \cap V \backslash \bigcup_{k \neq j} X_{k} \cap X_{j, \mathrm{reg}} \cap V
$$

which are dense in $X_{j, \text { reg }} \cap V$. Therefore any bounded holomorphic function on $X_{\text {reg }} \cap V$ extends to each component $X_{j, \text { reg }} \cap V$ and we see that

$$
\tilde{\mathscr{O}}_{X, x}=\bigoplus \tilde{\mathscr{O}}_{X_{j}, x} .
$$

The first important fact is that weakly holomorphic functions are always meromorphic and possess "universal denominators".
(7.2) Theorem. For every point $x \in X$, there is a neighborhood $V$ of $x$ and $h \in \mathbb{O}_{X}(V)$ such that $h^{-1}(0)$ is nowhere dense in $V$ and $h_{y} \tilde{\mathscr{O}}_{X, y} \subset \mathscr{O}_{X, y}$ for all $y \in V$; such a function $h$ is called a universal denominator on $V$. In particular $\tilde{\mathscr{O}}_{X}$ is contained in the ring $\mathscr{M}_{X}$ of meromorphic functions.

Proof. First assume that $(X, x)$ is irreducible and that we have a ramified covering $\pi: X \cap \Delta \longrightarrow \Delta^{\prime}$ with ramification locus $S$. We claim that the discriminant $\delta\left(z^{\prime}\right)$ of a
primitive element $u\left(z^{\prime \prime}\right)=c_{d+1} z_{d+1}+\cdots+c_{n} z_{n}$ is a universal denominator on $X \cap \Delta$. To see this, we imitate the proof of Lemma 4.15. Let $f \in \tilde{\mathscr{O}}_{X, y}, y \in X \cap \Delta$. Then we solve the equation

$$
f(z)=\sum_{0 \leqslant j \leqslant q} b_{j}\left(z^{\prime}\right) u\left(z^{\prime \prime}\right)^{j}
$$

in a neighborhood of $y$. For $z^{\prime} \in \Delta^{\prime} \backslash S$, let us denote by $\left(z^{\prime}, z_{\alpha}^{\prime \prime}\right), 1 \leqslant \alpha \leqslant q$, the points in the fiber $X \cap \pi^{-1}\left(z^{\prime}\right)$. Among these, only $q^{\prime}$ are close to $y$, where $q^{\prime}$ is the sum of the sheet numbers of the irreducible components of $(X, y)$ by the projection $\pi$. The other points $\left(z^{\prime}, z_{\alpha}^{\prime \prime}\right)$, say $q^{\prime}<\alpha \leqslant q$, are in neighborhoods of the points of $\pi^{-1}\left(y^{\prime}\right) \backslash\{y\}$. We take $\left(b_{j}\left(z^{\prime}\right)\right)$ to be the solution of the linear system

$$
\sum_{0 \leqslant j \leqslant q} b_{j}\left(z^{\prime}\right) u\left(z_{\alpha}^{\prime \prime}\right)^{j}= \begin{cases}f\left(z^{\prime}, z_{\alpha}^{\prime \prime}\right) & \text { for } 1 \leqslant \alpha \leqslant q^{\prime} \\ 0 & \text { for } q^{\prime}<\alpha \leqslant n\end{cases}
$$

The solutions $b_{j}\left(z^{\prime}\right)$ are holomorphic on $\Delta^{\prime} \backslash S$ near $y^{\prime}$. Since the determinant is $\delta\left(z^{\prime}\right)^{1 / 2}$, we see that $\delta b_{j}$ is bounded, thus $\delta b_{j} \in \mathbb{O}_{\Delta^{\prime}, y^{\prime}}$ and $\delta_{y} f \in \mathbb{O}_{X, y}$.

Now, assume that $(X, x) \subset\left(\mathbb{C}^{n}, 0\right)$ has irreducible components $\left(X_{j}, x\right)$. We can find for each $j$ a neighborhood $\Omega_{j}$ of 0 in $\mathbb{C}^{n}$ and a function $\delta_{j} \in \mathscr{O}_{n}\left(\Omega_{j}\right)$ which is a universal denominator on $X_{j} \cap \Omega_{j}$. After adding to $\delta_{j}$ a function which is identically zero on ( $X_{j}, x$ ) and non zero on ( $\left.X_{k}, x\right), k \neq j$, we may assume that $\delta_{j}^{-1}(0) \cap X_{k} \cap \Omega$ is nowhere dense in $X_{k} \cap \Omega$ for all $j$ and $k$ and some small $\Omega \subset \bigcap \Omega_{j}$. Then $\delta=\prod \delta_{j}$ is a universal denominator on each component $X_{j} \cap \Omega$. For some possibly smaller $\Omega$, select a function $v_{j} \in \mathscr{O}_{n}(\Omega)$ such that $v_{j}$ vanishes identically on $\bigcup_{k \neq j} X_{k} \cap \Omega$ and $v_{j}^{-1}(0)$ is nowhere dense in $X_{j} \cap \Omega$, and set $h=\delta \sum v_{k}$. For any germ $f \in \mathscr{O}_{X, y}, y \in X \cap \Omega$, there is a germ $g_{j} \in \mathfrak{Q}_{\Omega, y}$ with $\delta f=g_{j}$ on $\left(X_{j}, y\right)$. We have $h=\delta v_{j}$ on $X_{j} \cap \Omega$, so $h^{-1}(0)$ is nowhere dense in $X \cap \Omega$ and

$$
h f=v_{j} \delta f=v_{j} g_{j}=\sum v_{k} g_{k} \quad \text { on each } \quad\left(X_{j}, y\right) .
$$

Since $\sum v_{k} g_{k} \in \mathscr{O}_{\Omega, y}$, we get $h \tilde{\mathscr{O}}_{X, y} \subset \mathscr{O}_{X, y}$.
(7.3) Theorem. If $(X, x)$ is irreducible, $\tilde{\mathscr{O}}_{X, x}$ is the integral closure of $\mathscr{O}_{X, x}$ in its quotient field $\mathcal{M}_{X, x}$. Moreover, every germ $f \in \tilde{\mathscr{O}}_{X, x}$ admits a limit

$$
\lim _{X_{\operatorname{reg}} \ni z \rightarrow x} f(z) .
$$

Observe that $\mathscr{O}_{X, x}$ is an entire ring, so the ring of quotients $\mathscr{M}_{X, x}$ is actually a field. A simple illustration of the theorem is obtained with the irreducible germ of curve $X: z_{1}^{3}=z_{2}^{2}$ in $\left(\mathbb{C}^{2}, 0\right)$. Then $X$ can be parametrized by $z_{1}=t^{2}, z_{2}=t^{3}, t \in \mathbb{C}$, and $\mathfrak{O}_{X, 0}=\mathbb{C}\left\{z_{1}, z_{2}\right\} /\left(z_{1}^{3}-z_{2}^{2}\right)=\mathbb{C}\left\{t^{2}, t^{3}\right\}$ consists of all convergent series $\sum a_{n} t^{n}$ with $a_{1}=0$. The function $z_{2} / z_{1}=t$ is weakly holomorphic on $X$ and satisfies the integral equation $t^{2}-z_{1}=0$. Here we have $\tilde{\mathscr{O}}_{X, 0}=\mathbb{C}\{t\}$.

Proof. a) Let $f=g / h$ be an element in $\mathcal{M}_{X, x}$ satisfying an integral equation

$$
f^{m}+a_{1} f^{m-1}+\ldots+a_{m}=0, \quad a_{k} \in \mathbb{O}_{X, x} .
$$

Set $A=h^{-1}(0)$. Then $f$ is holomorphic on $X \backslash A$ near $x$, and Lemma 4.10 shows that $f$ is bounded on a neighborhood of $x$. By Remark $4.2, f$ can be extended as a holomorphic function on $X_{\text {reg }}$ in a neighborhood of $x$, thus $f \in \widetilde{\mathscr{O}}_{X, x}$.
b) Let $f \in \tilde{\mathscr{O}}_{X, x}$ and let $\pi: X \cap \Delta \longrightarrow \Delta^{\prime}$ be a ramified covering in a neighborhood of $x$, with ramification locus $S$. As in the proof of Th. 6.1, $f$ satisfies an equation

$$
f^{q}-\sigma_{1} f^{q-1}+\cdots+(-1)^{q} \sigma_{q}=0, \quad \sigma_{k} \in \mathbb{O}\left(\Delta^{\prime}\right) ;
$$

indeed the elementary symmetric functions $\sigma_{k}\left(z^{\prime}\right)$ are holomorphic on $\Delta^{\prime} \backslash S$ and bounded, so they extend holomorphically to $\Delta^{\prime}$. Hence $\tilde{\mathscr{O}}_{X, x}$ is integral over $\mathscr{O}_{X, x}$ and we already know that $\tilde{\mathscr{O}}_{X, x} \subset \mathcal{M}_{X, x}$.
c) Finally, the cluster set $\bigcap_{V \ni x} \overline{f\left(X_{\text {reg }} \cap V\right)}$ is connected, because there is a fundamental system of neighborhoods $V$ of $x$ such that $X_{\mathrm{reg}} \cap V$ is connected, and any intersection of a decreasing sequence of compact connected sets is connected. However the limit set is contained in the finite set of roots of equation b) at point $x^{\prime} \in \Delta^{\prime}$, so it must be reduced to one element.

## § 7.2. Normal Spaces

Normal spaces are spaces for which all weakly holomorphic functions are actually holomorphic. These spaces will be seen later to have "simpler" singularities than general analytic spaces.
(7.4) Definition. A complex space $X$ is said to be normal at a point $x$ if $(X, x)$ is irreducible and $\tilde{\mathscr{O}}_{X, x}=\mathscr{O}_{X, x}$, that is, $\mathscr{O}_{X, x}$ is integrally closed in its field of quotients. The set of normal (resp. non-normal) points will be denoted $X_{\mathrm{norm}}$ (resp. $X_{\mathrm{n}-\mathrm{n}}$ ). The space $X$ itself is said to be normal if $X$ is normal at every point.

Observe that any regular point $x$ is normal: in fact $\mathbb{O}_{X, x} \simeq \mathscr{O}_{n}$ is then factorial, hence integrally closed. Therefore $X_{\text {n-n }} \subset X_{\text {sing }}$.
(7.5) Theorem. The non-normal set $X_{\mathrm{n}-\mathrm{n}}$ is an analytic subset of $X$. In particular, $X_{\text {norm }}$ is open in $X$.

Proof. We give here a beautifully simple proof due to [Grauert and Remmert 1984]. Let $h$ be a universal denominator on a neighborhood $V$ of a given point and let $\mathscr{I}=\sqrt{h \mathbb{O}_{X}}$ be the sheaf of ideals of $h^{-1}(0)$ by Hilbert's Nullstellensatz. Finally, let $\mathscr{F}=\operatorname{hom}_{\mathscr{C}}(\mathscr{I}, \mathscr{I})$ be the sheaf of $\mathscr{O}_{X}$-endomorphisms of $\mathscr{I}$. Since $\mathscr{I}$ is coherent, so is $\mathscr{F}$ (cf. Exercise 10.?). Clearly, the homotheties of $\mathscr{I}$ give an injection $\mathbb{O}_{X} \subset \mathscr{F}$ over $V$. We claim that there is a natural injection $\mathscr{F} \subset \tilde{\mathscr{O}}_{X}$. In fact, any endomorphism of $\mathscr{I}$ yields by restriction a homomorphism $h \mathscr{O}_{X} \longrightarrow \mathscr{O}_{X}$, and by $\mathscr{O}_{X}$-linearity such a homomorphism is obtained by multiplication by an element in $h^{-1} \mathscr{O}_{X}$. Thus $\mathscr{F} \subset h^{-1} \mathscr{O}_{X} \subset \mathscr{M}_{X}$. Since each stalk $\mathscr{I}_{x}$ is a finite $\mathbb{0}_{X, x}$-module containing non-zero divisors, it follows that that any meromorphic germ $f$ such that $f \mathscr{I}_{x} \subset \mathscr{I}_{x}$ is integral over $\mathscr{O}_{X, x}$ ([Lang 1965], Chapter IX, § 1), hence $\mathscr{F}_{x} \subset \tilde{\mathscr{O}}_{X, x}$. Thus we have inclusions $\mathscr{O}_{X} \subset \mathscr{F} \subset \tilde{\mathscr{O}}_{X}$. Now, we assert that

$$
X_{\mathrm{n}-\mathrm{n}} \cap V=\left\{x \in V ; \mathscr{F}_{x} \neq \mathscr{O}_{X, x}\right\}=\mathscr{F} / \mathbb{O}_{X} .
$$

This will imply the theorem by 5.10. To prove the equality, we first observe that $\mathscr{F}_{x} \neq$ $\mathfrak{O}_{X, x}$ implies $\widetilde{\mathscr{O}}_{X, x} \neq \mathscr{O}_{X, x}$, thus $x \in X_{\mathrm{n}-\mathrm{n}}$. Conversely, assume that $x$ is non normal, that is, $\tilde{\mathscr{O}}_{X, x} \neq \mathfrak{O}_{X, x}$. Let $k$ be the smallest integer such that $\mathscr{I}_{x}^{k} \tilde{\mathscr{O}}_{X, x} \subset \mathfrak{O}_{X, x} ;$ such an integer exists since $\mathscr{\oiint}_{x}^{l} \tilde{\mathscr{O}}_{X, x} \subset h \tilde{\mathscr{O}}_{X, x} \subset \mathscr{O}_{X, x}$ for $l$ large. Then there is an element $w \in \mathscr{I}_{x}^{k-1} \tilde{\mathscr{O}}_{X, x}$ such that $w \notin \mathscr{O}_{X, x}$. We have $w \mathscr{I}_{x} \subset \mathscr{O}_{X, x} ;$ moreover, as $w$ is locally bounded near $X_{\text {sing }}$, any germ $w g$ in $w \mathscr{I}_{x}$ satisfies $\lim w(z) g(z)=0$ when $z \in X_{\text {reg }}$ tends to a point of the zero variety $h^{-1}(0)$ of $\mathscr{I}_{x}$. Hence $w \mathscr{I}_{x} \subset \mathscr{I}_{x}$, i.e. $w \in \mathscr{F}_{x}$, but $w \notin \mathbb{O}_{X, x}$, so $\mathscr{F}_{x} \neq \mathfrak{O}_{X, x}$.
(7.6) Theorem. If $x \in X$ is a normal point, then $\left(X_{\operatorname{sing}}, x\right)$ has codimension at least 2 in $(X, x)$.

Proof. We suppose that $\Sigma=X_{\text {sing }}$ has codimension 1 in a neighborhood of $x$ and try to get a contradiction. By restriction to a smaller neighborhood, we may assume that $X$ itself is normal and irreducible (since $X_{\text {norm }}$ is open), $\operatorname{dim} X=n$, that $\Sigma$ has pure dimension $n-1$ and that the ideal sheaf $\mathscr{I}_{\Sigma}$ has global generators $\left(g_{1}, \ldots, g_{k}\right)$. Then $\Sigma \subset \bigcup g_{j}^{-1}(0)$; both sets have pure dimension $n-1$ and thus singular sets of dimension $\leqslant n-2$. Hence there is a point $a \in \Sigma$ that is regular on $\Sigma$ and on $\bigcup g_{j}^{-1}(0)$, in particular there is a neighborhood $V$ of $a$ such that $g_{1}^{-1}(0) \cap V=\ldots=g_{k}^{-1}(0) \cap V=\Sigma \cap V$ is a smooth $(n-1)$-dimensional manifold. Since $\operatorname{codim}_{X} \Sigma=1$ and $a$ is a singular point of $X, \mathscr{I}_{\Sigma, a}$ cannot have less than 2 generators in $\mathbb{O}_{X, a}$ by Cor. 4.33. Take $\left(g_{1}, \ldots, g_{l}\right)$, $l \geqslant 2$, to be a minimal subset of generators. Then $f=g_{2} / g_{1}$ cannot belong to $\mathscr{O}_{X, a}$, but $f$ is holomorphic on $V \backslash \Sigma$. We may assume that there is a sequence $a_{\nu} \in V \backslash \Sigma$ converging to $a$ such that $f\left(a_{\nu}\right)$ remains bounded (otherwise reverse $g_{1}$ and $g_{2}$ and pass to a subsequence). Since $g_{1}^{-1}(0) \cap V=\Sigma \cap V$, Hilbert's Nullstellensatz gives an integer $m$ such that $\mathscr{I}_{\Sigma, a}^{m} \subset g_{1} \mathscr{O}_{X, a}$, hence $f_{a} \mathscr{I}_{\Sigma, a}^{m} \subset \mathscr{O}_{X, a}$. We take $m$ to be the smallest integer such that the latter inclusion holds. Then there is a product $g^{\alpha}=g_{1}^{\alpha_{1}} \ldots g_{l}^{\alpha_{l}}$ with $|\alpha|=m-1$ such that $f g^{\alpha} \notin \mathscr{O}_{X, a}$ but $f g^{\alpha} g_{j} \in \mathscr{O}_{X, a}$ for each $j$. Since the sequence $f\left(a_{\nu}\right)$ is bounded we conclude that $f g^{\alpha} g_{j}$ vanishes at $a$. The zero set of this function has dimension $n-1$ and is contained in $\bigcup g_{k}^{-1}(0) \cap V=\Sigma \cap V$ so it must contain the germ ( $\Sigma, a)$. Hence $f g^{\alpha} g_{j} \in \mathscr{I}_{\Sigma, a}$ and $f g^{\alpha} \mathscr{I}_{\Sigma, a} \subset \mathscr{I}_{\Sigma, a}$. As $\mathscr{I}_{\Sigma, a}$ is a finitely generated $\mathfrak{O}_{X, a}$-module, this implies $f g^{\alpha} \in \tilde{\mathscr{O}}_{X, a}=\mathscr{O}_{X, a}$, a contradiction.
(7.7) Corollary. A complex curve is normal if and only if it is regular.
(7.8) Corollary. Let $X$ be a normal complex space and $Y$ an analytic subset of $X$ such that $\operatorname{dim}(Y, x) \leqslant \operatorname{dim}(X, x)-2$ for any $x \in X$. Then any holomorphic function on $X \backslash Y$ can be extended to a holomorphic function on $X$.

Proof. By Cor. 1.4.5, every holomorphic function $f$ on $X_{\text {reg }} \backslash Y$ extends to $X_{\text {reg }}$. Since codim $X_{\text {sing }} \geqslant 2$, Th. 6.1 shows that $f$ is locally bounded near $X_{\text {sing }}$. Therefore $f$ extends to $X$ by definition of a normal space.

## § 7.3. The Oka Normalization Theorem

The important normalization theorem of [Oka 1950] shows that $\tilde{\mathscr{O}}_{X}$ can be used to define the structure sheaf of a new analytic space $\tilde{X}$ which is normal and is obtained by "simplifying" the singular set of $X$. More precisely:
(7.9) Definition. Let $X$ be a complex space. A normalization $(Y, \pi)$ of $X$ is a normal complex space $Y$ together with a holomorphic map $\pi: Y \longrightarrow X$ such that the following conditions are satisfied.
a) $\pi: Y \longrightarrow X$ is proper and has finite fibers;
b) if $\Sigma$ is the set of singular points of $X$ and $A=\pi^{-1}(\Sigma)$, then $Y \backslash A$ is dense in $Y$ and $\pi: Y \backslash A \longrightarrow X \backslash \Sigma=X_{\text {reg }}$ is an analytic isomorphism.

It follows from b) that $Y \backslash A \subset Y_{\text {reg }}$. Thus $Y$ is obtained from $X$ by a suitable "modification" of its singular points. Observe that $Y_{\text {reg }}$ may be larger than $Y \backslash A$, as is the case in the following two examples.

## (7.10) Examples.

a) Let $X=\mathbb{C} \times\{0\} \cup\{0\} \times \mathbb{C}$ be the complex curve $z_{1} z_{2}=0$ in $\mathbb{C}^{2}$. Then the normalization of $X$ is the disjoint union $Y=\mathbb{C} \times\{1,2\}$ of two copies of $\mathbb{C}$, with the map $\pi\left(t_{1}\right)=\left(t_{1}, 0\right)$, $\pi\left(t_{2}\right)=\left(0, t_{2}\right)$. The set $A=\pi^{-1}(0,0)$ consists of exactly two points.
b) The cubic curve $X: z_{1}^{3}=z_{2}^{2}$ is normalized by the map $\pi: \mathbb{C} \longrightarrow X, t \longmapsto\left(t^{2}, t^{3}\right)$. Here $\pi$ is a homeomorphism but $\pi^{-1}$ is not analytic at $(0,0)$.

We first show that the normalization is essentially unique up to isomorphism and postpone the proof of its existence for a while.
(7.11) Lemma. If $\left(Y_{1}, \pi_{1}\right)$ and $\left(Y_{2}, \pi_{2}\right)$ are normalizations of $X$, there is a unique analytic isomorphism $\varphi: Y_{1} \longrightarrow Y_{2}$ such that $\pi_{1}=\pi_{2} \circ \varphi$.

Proof. Let $\Sigma$ be the set of singular points of $X$ and $A_{j}=\pi_{j}^{-1}(\Sigma), j=1,2$. Let $\varphi^{\prime}: Y_{1} \backslash A_{1} \longrightarrow Y_{2} \backslash A_{2}$ be the analytic isomorphism $\pi_{2}^{-1} \circ \pi_{1}$. We assert that $\varphi^{\prime}$ can be extended into a map $\varphi: Y_{1} \longrightarrow Y_{2}$. In fact, let $a \in A_{1}$ and $s=\pi_{1}(a) \in \Sigma$. Then $\pi_{2}^{-1}(s)$ consists of a finite set of points $y_{j} \in Y_{2}$. Take disjoint neighborhoods $U_{j}$ of $y_{j}$ such that $U_{j}$ is an analytic subset in an open set $\Omega_{j} \subset \subset \mathbb{C}^{N}$. Since $\pi_{2}$ is proper, there is a neighborhood $V$ of $s$ in $X$ such that $\pi_{2}^{-1}(V) \subset \bigcup U_{j}$ and by continuity of $\pi_{1}$ a neighborhood $W$ of $a$ such that $\pi_{1}(W) \subset V$. Then $\varphi^{\prime}=\pi_{2}^{-1} \circ \pi_{1}$ maps $W \backslash A_{1}$ into $\bigcup U_{j}$ and can be seen as a bounded holomorphic map into $\mathbb{C}^{N}$ through the embeddings $U_{j} \subset \Omega_{j} \subset \subset \mathbb{C}^{N}$. Since $Y_{1}$ is normal, $\varphi^{\prime}$ extends to $W$, and the extension takes values in $\bigcup \bar{U}_{j}$ which is contained in $Y_{2}$ (shrink $U_{j}$ if necessary). Thus $\varphi^{\prime}$ extends into a map $\varphi: Y_{1} \longrightarrow Y_{2}$ and similarly $\varphi^{\prime-1}$ extends into a map $\psi: Y_{2} \longrightarrow Y_{1}$. By density of $Y_{j} \backslash A_{j}$, we have $\psi \circ \varphi=\operatorname{Id}_{Y_{1}}, \varphi \circ \psi=\operatorname{Id}_{Y_{2}}$.
(7.12) Oka normalization theorem. Let $X$ be any complex space. Then $X$ has a normalization $(Y, \pi)$.

Proof. Because of the previous lemma, it suffices to prove that any point $x \in X$ has a neighborhood $U$ such that $U$ admits a normalization; all these local normalizations will then glue together. Hence we may suppose that $X$ is an analytic set in an open set of $\mathbb{C}^{n}$. Moreover, if $(X, x)$ splits into irreducible components $\left(X_{j}, x\right)$ and if $\left(Y_{j}, \pi_{j}\right)$ is a normalization of $X_{j} \cap U$, then the disjoint union $Y=\coprod Y_{j}$ with $\pi=\coprod \pi_{j}$ is easily seen to be a normalization of $X \cap U$. We may therefore assume that $(X, x)$ is irreducible.

Let $h$ be a universal denominator in a neighborhood of $x$. Then $\tilde{\mathscr{O}}_{X, x}$ is isomorphic to its image $h \tilde{\mathscr{O}}_{X, x} \subset \mathfrak{O}_{X, x}$, so it is a finitely generated $\mathbb{O}_{X, x}$-module. Let $\left(f_{1}, \ldots, f_{m}\right)$ be a finite set of generators of $\mathbb{O}_{X, x}$. After shrinking $X$ again, we may assume the following two points:

- $X$ is an analytic set in an open set $\Omega \subset \mathbb{C}^{n},(X, x)$ is irreducible and $X_{\text {reg }}$ is connected;
- $f_{j}$ is holomorphic in $X_{\text {reg }}$, can be written $f_{j}=g_{j} / h$ on $X$ with $g_{j}, h$ in $\mathscr{O}_{n}(\Omega)$ and satisfies an integral equation $P_{j}\left(z ; f_{j}(z)\right)=0$ where $P_{j}(z ; T)$ is a unitary polynomial with holomorphic coefficients on $X$.

Set $X^{\prime}=X \backslash h^{-1}(0)$. Consider the holomorphic map

$$
F: X_{\mathrm{reg}} \longrightarrow \Omega \times \mathbb{C}^{m}, \quad z \longmapsto\left(z, f_{1}(z), \ldots, f_{m}(z)\right)
$$

and the image $Y^{\prime}=F\left(X^{\prime}\right)$. We claim that the closure $Y$ of $Y^{\prime}$ in $\Omega \times \mathbb{C}^{m}$ is an analytic set. In fact, the set

$$
Z=\left\{(z, w) \in \Omega \times \mathbb{C}^{m} ; z \in X, h(z) w_{j}=g_{j}(z)\right\}
$$

is analytic and $Y^{\prime}=Z \backslash\{h(z)=0\}$, so we may apply Cor. 5.4. Observe that $Y^{\prime}$ is contained in the set defined by $P_{j}\left(z ; w_{j}\right)=0$, thus so is its closure $Y$. The first projection $\Omega \times \mathbb{C}^{m} \longrightarrow \Omega$ gives a holomorphic map $\pi: Y \longrightarrow X$ such that $\pi \circ F=\operatorname{Id}$ on $X^{\prime}$, hence also on $X_{\text {reg }}$. If $\Sigma=X_{\text {sing }}$ and $A=\pi^{-1}(\Sigma)$, the restriction $\pi: Y \backslash A \longrightarrow X \backslash \Sigma=X_{\text {reg }}$ is thus an analytic isomorphism and $F$ is its inverse. Since $(X, x)$ is irreducible, each $f_{j}$ has a limit $\ell_{j}$ at $x$ by Th. 7.3 and the fiber $\pi^{-1}(x)$ is reduced to the single point $y=(x, \ell)$. The other fibers $\pi^{-1}(z)$ are finite because they are contained in the finite set of roots of the equations $P_{j}\left(z ; w_{j}\right)=0$. The same argument easily shows that $\pi$ is proper (use Lemma 4.10).

Next, we show that $Y$ is normal at the point $y=\pi^{-1}(x)$. In fact, for any bounded holomorphic function $u$ on $\left(Y_{\text {reg }}, y\right)$ the function $u \circ F$ is bounded and holomorphic on $\left(X_{\text {reg }}, x\right)$. Hence $u \circ F \in \widetilde{\mathscr{O}}_{X, x}=\mathscr{O}_{X, x}\left[f_{1}, \ldots, f_{m}\right]$ and we can write $u \circ F(z)=$ $Q\left(z ; f_{1}(z), \ldots, f_{m}(z)\right)=Q \circ F(z)$ where $Q(z ; w)=\sum a_{\alpha}(z) w^{\alpha}$ is a polynomial in $w$ with coefficients in $\mathbb{0}_{X, x}$. Thus $u$ coincides with $Q$ on $\left(Y_{\mathrm{reg}}, y\right)$, and as $Q$ is holomorphic on $(X, x) \times \mathbb{C}^{m} \supset(Y, y)$, we conclude that $u \in \mathscr{O}_{Y, y}$. Therefore $\tilde{\mathscr{O}}_{Y, y}=\mathscr{O}_{Y, y}$.

Finally, by Th. 7.5, there is a neighborhood $V \subset Y$ of $y$ such that every point of $V$ is normal. As $\pi$ is proper, we can find a neighborhood $U$ of $x$ with $\pi^{-1}(U) \subset V$. Then $\pi: \pi^{-1}(U) \longrightarrow U$ is the required normalization in a neighborhood of $x$.

The proof of Th. 7.12 shows that the fiber $\pi^{-1}(x)$ has exactly one point $y_{j}$ for each irreducible component $\left(X_{j}, x\right)$ of $(X, x)$. As a one-to-one proper map is a homeomorphism, we get in particular:
(7.13) Corollary. If $X$ is a locally irreducible complex space, the normalization $\pi$ : $Y \longrightarrow X$ is a homeomorphism.
(7.14) Remark. In general, for any open set $U \subset X$, we have an isomorphism

$$
\begin{equation*}
\pi^{\star}: \tilde{\mathscr{O}}_{X}(U) \xrightarrow{\simeq} \mathscr{O}_{Y}\left(\pi^{-1}(U)\right), \tag{7.15}
\end{equation*}
$$

whose inverse is given by the comorphism of $\pi^{-1}: X_{\text {reg }} \longrightarrow Y$; note that $\tilde{\mathscr{O}}_{Y}(U)=\mathscr{O}_{Y}(U)$ since $Y$ is normal. Taking the direct limit over all neighborhoods $U$ of a given point $x \in X$, we get an isomorphism

$$
\begin{equation*}
\pi^{\star}: \tilde{\mathscr{O}}_{X, x} \longrightarrow \bigoplus_{y_{j} \in \pi^{-1}(x)} \mathscr{O}_{Y, y_{j}} \tag{7.15'}
\end{equation*}
$$

In other words, $\tilde{\mathscr{O}}_{X}$ is isomorphic to the direct image sheaf $\pi_{\star} \mathscr{O}_{Y}$, see (1.12). We will prove later on the deep fact that the direct image of a coherent sheaf by a proper holomorphic map is always coherent ([Grauert 1960], see 9.?.1). Hence $\tilde{\mathscr{O}}_{X}=\pi_{\star} \mathscr{O}_{Y}$ is a coherent sheaf over $\mathbb{O}_{X}$.

## § 8. Holomorphic Mappings and Extension Theorems

## § 8.1. Rank of a Holomorphic Mapping

Our goal here is to introduce the general concept of the rank of a holomorphic map and to relate the rank to the dimension of the fibers. As in the smooth case, the rank is shown to satisfy semi-continuity properties.
(8.1) Lemma. Let $F: X \longrightarrow Y$ be a holomorphic map from a complex space $X$ to $a$ complex space $Y$.
a) If $F$ is finite, i.e. proper with finite fibers, then $\operatorname{dim} X \leqslant \operatorname{dim} Y$.
b) If $F$ is finite and surjective, then $\operatorname{dim} X=\operatorname{dim} Y$.

Proof. a) Let $x \in X,\left(X_{j}, x\right)$ an irreducible component and $m=\operatorname{dim}\left(X_{j}, x\right)$. If $\left(Y_{k}, y\right)$ are the irreducible components of $Y$ at $y=F(x)$, then $\left(X_{j}, x\right)$ is contained in $\bigcup F^{-1}\left(Y_{k}\right)$, hence $\left(X_{j}, x\right)$ is contained in one of the sets $F^{-1}\left(Y_{k}\right)$. If $p=\operatorname{dim}\left(Y_{k}, y\right)$, there is a ramified covering $\pi$ from some neighborhood of $y$ in $Y_{k}$ onto a polydisk in $\Delta^{\prime} \subset \mathbb{C}^{p}$. Replacing $X$ by some neighborhood of $x$ in $X_{j}$ and $F$ by the finite map $\pi \circ F_{\left\lceil X_{j}\right.}: X_{j} \longrightarrow \Delta^{\prime}$, we may suppose that $Y=\Delta^{\prime}$ and that $X$ is irreducible, $\operatorname{dim} X=m$. Let $r=\operatorname{rank} d F_{x_{0}}$ be the maximum of the rank of the differential of $F$ on $X_{\text {reg }}$. Then $r \leqslant \min \{m, p\}$ and the rank of $d F$ is constant equal to $r$ on a neighborhood $U$ of $x_{0}$. The constant rank theorem implies that the fibers $F^{-1}(y) \cap U$ are ( $m-r$ )-dimensional submanifolds, hence $m-r=0$ and $m=r \leqslant p$.
b) We only have to show that $\operatorname{dim} X \geqslant \operatorname{dim} Y$. Fix a regular point $y \in Y$ of maximal dimension. By taking the restriction $F: F^{-1}(U) \longrightarrow U$ to a small neighborhood $U$ of $y$, we may assume that $Y$ is an open subset of $\mathbb{C}^{p}$. If $\operatorname{dim} X<\operatorname{dim} Y$, then $X$ is a union of analytic manifolds of dimension $<\operatorname{dim} Y$ and Sard's theorem implies that $F(X)$ has zero Lebesgue measure in $Y$, a contradiction.
(8.2) Proposition. For any holomorphic map $F: X \longrightarrow Y$, the fiber dimension $\operatorname{dim}\left(F^{-1}(F(x)), x\right)$ is an upper semi-continuous function of $x$.

Proof. Without loss of generality, we may suppose that $X$ is an analytic set in $\Omega \subset \mathbb{C}^{n}$, that $F(X)$ is contained in a small neighborhood of $F(x)$ in $Y$ which is embedded in $\mathbb{C}^{N}$,
and that $x=0, F(x)=0$. Set $A=F^{-1}(0)$ and $s=\operatorname{dim}(A, 0)$. We can find a linear form $\xi_{1}$ on $\mathbb{C}^{n}$ such that $\operatorname{dim}\left(A \cap \xi_{1}^{-1}(0), 0\right)=s-1$; in fact we need only select a point $x_{j} \neq 0$ on each irreducible component $\left(A_{j}, 0\right)$ of $(A, 0)$ and take $\xi_{1}\left(x_{j}\right) \neq 0$. By induction, we can find linearly independent forms $\xi_{1}, \ldots, \xi_{s}$ on $\mathbb{C}^{n}$ such that

$$
\operatorname{dim}\left(A \cap \xi_{1}^{-1}(0) \cap \ldots \cap \xi_{j}^{-1}(0), 0\right)=s-j
$$

for all $j=1, \ldots, s$; in particular 0 is an isolated point in the intersection when $j=s$. After a change of coordinates, we may suppose that $\xi_{j}(z)=z_{j}$. Fix $r^{\prime \prime}$ so small that the ball $\bar{B}^{\prime \prime} \subset \mathbb{C}^{n-s}$ of center 0 and radius $r^{\prime \prime}$ satisfies $A \cap\left(\{0\} \times \bar{B}^{\prime \prime}\right)=\{0\}$. Then $A$ is disjoint from the compact set $\{0\} \times \partial B^{\prime \prime}$, so there exists a small ball $B^{\prime} \subset \mathbb{C}^{s}$ of center 0 such that $A \cap\left(\bar{B}^{\prime} \times \partial B^{\prime \prime}\right)=\emptyset$, i.e. $F$ does not vanish on the compact set $K=X \cap\left(\bar{B}^{\prime} \times \partial B^{\prime \prime}\right)$. Set $\varepsilon=\min _{K}|F|$. Then for $|y|<\varepsilon$ the fiber $F^{-1}(y)$ does not intersect $\bar{B}^{\prime} \times \partial B^{\prime \prime}$. This implies that the projection map $\pi: F^{-1}(y) \cap\left(B^{\prime} \times B^{\prime \prime}\right) \longrightarrow B^{\prime}$ is proper. The fibers of $\pi$ are then compact analytic subsets of $B^{\prime \prime}$, so they are finite by 5.9. Lemma 8.1 a) implies

$$
\operatorname{dim} F^{-1}(y) \cap\left(B^{\prime} \times B^{\prime \prime}\right) \leqslant \operatorname{dim} B^{\prime}=s=\operatorname{dim}(A, 0)=\operatorname{dim}\left(F^{-1}(0), 0\right)
$$

Let $X$ be a pure dimensional complex space and $F: X \longrightarrow Y$ a holomorphic map. For any point $x \in X$, we define the rank of $F$ at $x$ by

$$
\begin{equation*}
\rho_{F}(x)=\operatorname{dim}(X, x)-\operatorname{dim}\left(F^{-1}(F(x)), x\right) . \tag{8.3}
\end{equation*}
$$

By the above proposition, $\rho_{F}$ is a lower semi-continuous function on $X$. In particular, if $\rho_{F}$ is maximum at some point $x_{0}$, it must be constant in a neighborhood of $x_{0}$. The maximum $\bar{\rho}(F)=\max _{X} \rho_{F}$ is thus attained on $X_{\text {reg }}$ or on any dense open subset $X^{\prime} \subset$ $X_{\text {reg }}$. If $X$ is not pure dimensional, we define $\bar{\rho}(F)=\max _{\alpha} \bar{\rho}\left(F_{\uparrow X_{\alpha}}\right)$ where $\left(X_{\alpha}\right)$ are the irreducible components of $X$. For a map $F: X \longrightarrow \mathbb{C}^{N}$, the constant rank theorem implies that $\bar{\rho}(F)$ is equal to the maximum of the rank of the jacobian matrix $d F$ at points of $X_{\text {reg }}$ (or of $X^{\prime}$ ).
(8.4) Proposition. If $F: X \longrightarrow Y$ is a holomorphic map and $Z$ an analytic subset of $X$, then $\bar{\rho}\left(F_{\upharpoonright Z}\right) \leqslant \bar{\rho}(F)$.

Proof. Since each irreducible component of $Z$ is contained in an irreducible component of $X$, we may assume $X$ irreducible. Let $\pi: \tilde{X} \longrightarrow X$ be the normalization of $X$ and $\tilde{Z}=\pi^{-1}(Z)$. Since $\pi$ is finite and surjective, the fiber of $F \circ \pi$ at point $x$ has the same dimension than the fiber of $F$ at $\pi(x)$ by Lemma 8.1 b$)$. Therefore $\bar{\rho}(F \circ \pi)=\bar{\rho}(F)$ and $\bar{\rho}\left(F \circ \pi_{\upharpoonright \tilde{Z}}\right)=\bar{\rho}\left(F_{\upharpoonright Z}\right)$, so we may assume $X$ normal. By induction on $\operatorname{dim} X$, we may also suppose that $Z$ has pure codimension 1 in $X$ (every point of $Z$ has a neighborhood $V \subset X$ such that $Z \cap V$ is contained in a pure one codimensional analytic subset of $V$ ). But then $Z_{\text {reg }} \cap X_{\text {reg }}$ is dense in $Z_{\text {reg }}$ because codim $X_{\text {sing }} \geqslant 2$. Thus we are reduced to the case when $X$ is a manifold and $Z$ a submanifold, and this case is clear if we consider the rank of the jacobian matrix.
(8.5) Theorem. Let $F: X \longrightarrow Y$ be a holomorphic map. If $Y$ is pure dimensional and $\bar{\rho}(F)<\operatorname{dim} Y$, then $F(X)$ has empty interior in $Y$.

Proof. Taking the restriction of $F$ to $F^{-1}\left(Y_{\text {reg }}\right)$, we may assume that $Y$ is a manifold. Since $X$ is a countable union of compact sets, so is $F(X)$, and Baire's theorem shows that the result is local for $X$. By Prop. 8.4 and an induction on $\operatorname{dim} X, F\left(X_{\text {sing }}\right)$ has empty interior in $Y$. The set $Z \subset X_{\text {reg }}$ of points where the jacobian matrix of $F$ has rank $<\bar{\rho}(F)$ is an analytic subset hence, by induction again, $F(Z)$ has empty interior. The constant rank theorem finally shows that every point $x \in X_{\text {reg }} \backslash Z$ has a neighborhood $V$ such that $F(V)$ is a submanifold of dimension $\bar{\rho}(F)$ in $Y$, thus $F(V)$ has empty interior and Baire's theorem completes the proof.
(8.6) Corollary. Let $F: X \longrightarrow Y$ be a surjective holomorphic map. Then $\operatorname{dim} Y=$ $\bar{\rho}(F)$.

Proof. By the remark before Prop. 8.4, there is a regular point $x_{0} \in X$ such that the jacobian matrix of $F$ has rank $\bar{\rho}(F)$. Hence, by the constant rank theorem $\operatorname{dim} Y \geqslant \bar{\rho}(F)$. Conversely, let $Y_{\alpha}$ be an irreducible component of $Y$ of dimension equal to $\operatorname{dim} Y$, and $Z=F^{-1}\left(Y_{\alpha}\right) \subset X$. Then $F(Z)=Y_{\alpha}$ and Th. 8.5 implies $\bar{\rho}(F) \geqslant \bar{\rho}\left(F_{\lceil Z}\right) \geqslant \operatorname{dim} Y_{\alpha}$.

## § 8.2. Remmert and Remmert-Stein Theorems

We are now ready to prove two important results: the extension theorem for analytic subsets due to [Remmert and Stein 1953] and the theorem of [Remmert 1956, 1957] which asserts that the image of a complex space under a proper holomorphic map is an analytic set. These will be obtained by a simultaneous induction on the dimension.
(8.7) Remmert-Stein theorem. Let $X$ be a complex space, $A$ an analytic subset of $X$ and $Z$ an analytic subset of $X \backslash A$. Suppose that there is an integer $p \geqslant 0$ such that $\operatorname{dim} A \leqslant p$, while $\operatorname{dim}(Z, x)>p$ for all $x \in Z$. Then the closure $\bar{Z}$ of $Z$ in $X$ is an analytic subset.
(8.8) Remmert's proper mapping theorem. Let $F: X \longrightarrow Y$ be a proper holomorphic map. Then $F(X)$ is an analytic subset of $Y$.

Proof. We let ( $8.7_{m}$ ) denote statement (8.7) for $\operatorname{dim} Z \leqslant m$ and ( $8.8_{m}$ ) denote statement (8.8) for $\operatorname{dim} X \leqslant m$. We proceed by induction on $m$ in two steps:

Step 1. $\left(8.7_{m}\right)$ and ( $8.8_{m-1}$ ) imply ( $8.8_{m}$ ).
Step 2. $\left(8.8_{m-1}\right)$ implies $\left(8.7_{m}\right)$.
As $\left(8.8_{m}\right)$ is obvious for $m=0$, our statements will then be valid for all $m$, i.e. for all complex spaces of bounded dimension. However, Th. 8.7 is local on $X$ and Th. 8.8 is local on $Y$, so the general case is immediately reduced to the finite dimensional case.

Proof of step 1. The analyticity of $F(X)$ is a local question in $Y$. Since $F: F^{-1}(U) \longrightarrow U$ is proper for any open set $U \subset Y$ and $F^{-1}(U) \subset \subset X$ if $U \subset \subset Y$, we may suppose that $Y$ is embedded in an open set $\Omega \subset \mathbb{C}^{n}$ and that $X$ only has finitely many irreducible components $X_{\alpha}$. Then we have $F(X)=\bigcup F\left(X_{\alpha}\right)$ and we are reduced to the case when $X$ is irreducible, $\operatorname{dim} X=m$ and $Y=\Omega$.

First assume that $X$ is a manifold and that the rank of $d F$ is constant. The constant rank theorem implies that every point in $X$ has a neighborhood $V$ such that $F(V)$ is a closed submanifold in a neighborhood $W$ of $F(x)$ in $Y$. For any point $y \in Y$, the
fiber $F^{-1}(y)$ can be covered by finitely many neighborhoods $V_{j}$ of points $x_{j} \in F^{-1}(y)$ such that $F\left(V_{j}\right)$ is a closed submanifold in a neighborhood $W_{j}$ of $y$. Then there is a neighborhood of $y W \subset \bigcap W_{j}$ such that $F^{-1}(W) \subset \bigcup V_{j}$, so $F(X) \cap W=\bigcup F\left(V_{j}\right) \cap W$ is a finite union of closed submanifolds in $W$ and $F(X)$ is analytic in $Y$.

Now suppose that $X$ is a manifold, set $r=\bar{\rho}(F)$ and let $Z \subset X$ be the analytic subset of points $x$ where the rank of $d F_{x}$ is $<r$. Since $\operatorname{dim} Z<m=\operatorname{dim} X$, the hypothesis $\left(8.8_{m-1}\right)$ shows that $F(Z)$ is analytic. We have $\operatorname{dim} F(Z)=\bar{\rho}\left(F_{\upharpoonright Z}\right)<r$. If $F(Z)=$ $F(X)$, then $F(X)$ is analytic. Otherwise $A=F^{-1}(F(Z))$ is a proper analytic subset of $X, d F$ has constant rank on $X \backslash A \subset X \backslash Z$ and the morphism $F: X \backslash A \longrightarrow Y \backslash F(Z)$ is proper. Hence the image $F(X \backslash A)$ is analytic in $Y \backslash F(Z)$. Since $\operatorname{dim} F(X \backslash A)=r \leqslant m$ and $\operatorname{dim} F(Z)<r$, hypothesis $\left(8.7_{m}\right)$ implies that $F(X)=\overline{F(X \backslash A)}$ is analytic in $Y$. When $X$ is not a manifold, we apply the same reasoning with $Z=X_{\text {sing }}$ in order to be reduced to the case of $F: X \backslash A \longrightarrow Y \backslash F(Z)$ where $X \backslash A$ is a manifold.

Proof of step 2. Since Th. 8.7 is local on $X$, we may suppose that $X$ is an open set $\Omega \subset \mathbb{C}^{n}$. Then we use induction on $p$ to reduce the situation to the case when $A$ is a $p$-dimensional submanifold (if this case is taken for granted, the closure of $Z$ in $\Omega \backslash A_{\text {sing }}$ is analytic and we conclude by the induction hypothesis). By a local analytic change of coordinates, we may assume that $0 \in A$ and that $A=\Omega \cap L$ where $L$ is a vector subspace of $\mathbb{C}^{n}$ of dimension $p$. By writing $Z=\bigcup_{p<s \leqslant m} Z_{s}$ where $Z_{s}$ is an analytic subset of $\Omega \backslash Y$ of pure dimension $s$, we may suppose that $Z$ has pure dimension $s, p<s \leqslant m$. We are going to show that $\bar{Z}$ is analytic in a neighborhood of 0 .

Let $\xi_{1}$ be a linear form on $\mathbb{C}^{n}$ which is not identically zero on $L$ nor on any irreducible component of $Z$ (just pick a point $x_{\nu}$ on each component and take $\xi_{1}\left(x_{\nu}\right) \neq 0$ for all $\nu$ ). Then $\operatorname{dim} L \cap \xi_{1}^{-1}(0)=p-1$ and the analytic set $Z \cap \xi_{1}^{-1}(0)$ has pure dimension $s-1$. By induction, there exist linearly independent forms $\xi_{1}, \ldots, \xi_{s}$ such that

$$
\begin{array}{cl}
\operatorname{dim} L \cap \xi_{1}^{-1}(0) \cap \ldots \cap \xi_{j}^{-1}(0)=p-j, & 1 \leqslant j \leqslant p \\
\operatorname{dim} Z \cap \xi_{1}^{-1}(0) \cap \ldots \cap \xi_{j}^{-1}(0)=s-j, & 1 \leqslant j \leqslant s \tag{8.9}
\end{array}
$$

By adding a suitable linear combination of $\xi_{1}, \ldots, \xi_{p}$ to each $\xi_{j}, p<j \leqslant s$, we may take $\xi_{j \upharpoonright L}=0$ for $p<j \leqslant s$. After a linear change of coordinates, we may suppose that $\xi_{j}(z)=z_{j}, L=\mathbb{C}^{p} \times\{0\}$ and $A=\Omega \cap\left(\mathbb{C}^{p} \times\{0\}\right)$. Let $\xi=\left(\xi_{1}, \ldots, \xi_{s}\right): \mathbb{C}^{n} \longrightarrow \mathbb{C}^{s}$ be the projection onto the first $s$ variables. As $Z$ is closed in $\Omega \backslash A, Z \cup A$ is closed in $\Omega$. Moreover, our construction gives $(Z \cup A) \cap \xi^{-1}(0)=\left(Z \cap \xi^{-1}(0)\right) \cup\{0\}$ and the case $j=s$ of (8.9) shows that $Z \cap \xi^{-1}(0)$ is a locally finite sequence in $\Omega \cap\left(\{0\} \times \mathbb{C}^{n-s}\right) \backslash\{0\}$. Therefore, we can find a small ball $\bar{B}^{\prime \prime}$ of center 0 in $\mathbb{C}^{n-s}$ such that $Z \cap\left(\{0\} \times \partial B^{\prime \prime}\right)=\emptyset$. As $\{0\} \times \partial B^{\prime \prime}$ is compact and disjoint from the closed set $Z \cup A$, there is a small ball $B^{\prime}$ of center 0 in $\mathbb{C}^{s}$ such that $(Z \cup A) \cap\left(\bar{B}^{\prime} \times \partial B^{\prime \prime}\right)=\emptyset$. This implies that the projection $\xi:(Z \cup A) \cap\left(B^{\prime} \times B^{\prime \prime}\right) \longrightarrow B^{\prime}$ is proper. Set $A^{\prime}=B^{\prime} \cap\left(\mathbb{C}^{p} \times\{0\}\right)$.


Fig. II-3 Projection $\pi: Z \cap\left(\left(B^{\prime} \backslash A^{\prime}\right) \times B^{\prime \prime}\right) \longrightarrow B^{\prime} \backslash A^{\prime}$.

Then the restriction

$$
\pi=\xi: Z \cap\left(B^{\prime} \times B^{\prime \prime}\right) \backslash\left(A^{\prime} \times B^{\prime \prime}\right) \longrightarrow B^{\prime} \backslash A^{\prime}
$$

is proper, and $Z \cap\left(B^{\prime} \times B^{\prime \prime}\right)$ is analytic in $\left(B^{\prime} \times B^{\prime \prime}\right) \backslash A$, so $\pi$ has finite fibers by Th. 5.9. By definition of the rank we have $\bar{\rho}(\pi)=s$. Let $S_{1}=Z_{\text {sing }} \cap \pi^{-1}\left(B^{\prime} \backslash A^{\prime}\right)$ and $S_{1}^{\prime}=\pi\left(S_{1}\right)$; further, let $S_{2}$ be the set of points $x \in Z \cap \pi^{-1}\left(B^{\prime} \backslash\left(A^{\prime} \cup S_{1}^{\prime}\right)\right) \subset Z_{\mathrm{reg}}$ such that $d \pi_{x}$ has rank $<s$ and $S_{2}^{\prime}=\pi\left(S_{2}\right)$. We have $\operatorname{dim} S_{j} \leqslant s-1 \leqslant m-1$. Hypothesis $(8.8)_{m-1}$ implies that $S_{1}^{\prime}$ is analytic in $B^{\prime} \backslash A^{\prime}$ and that $S_{2}^{\prime}$ is analytic in $B^{\prime} \backslash\left(A^{\prime} \cup S_{1}^{\prime}\right)$. By Remark 4.2, $B^{\prime} \backslash\left(A^{\prime} \cup S_{1}^{\prime} \cup S_{2}^{\prime}\right)$ is connected and every bounded holomorphic function on this set extends to $B^{\prime}$. As $\pi$ is a (non ramified) covering over $B^{\prime} \backslash\left(A^{\prime} \cup S_{1}^{\prime} \cup S_{2}^{\prime}\right)$, the sheet number is a constant $q$.

Let $\lambda(z)=\sum_{j>s} \lambda_{j} z_{j}$ be a linear form on $\mathbb{C}^{n}$ in the coordinates of index $j>s$. For $z^{\prime} \in B^{\prime} \backslash\left(A^{\prime} \cup S_{1}^{\prime} \cup S_{2}^{\prime}\right)$, we let $\sigma_{j}\left(z^{\prime}\right)$ be the elementary symmetric functions in the $q$ complex numbers $\lambda(z)$ corresponding to $z \in \pi^{-1}\left(z^{\prime}\right)$. Then these functions can be extended as bounded holomorphic functions on $B^{\prime}$ and we get a polynomial $P_{\lambda}\left(z^{\prime} ; T\right)$ such that $P_{\lambda}\left(z^{\prime} ; \lambda\left(z^{\prime \prime}\right)\right)$ vanishes identically on $Z \backslash \pi^{-1}\left(A^{\prime} \cup S_{1}^{\prime} \cup S_{2}^{\prime}\right)$. Since $\pi$ is finite, $Z \cap \pi^{-1}\left(A^{\prime} \cup S_{1}^{\prime} \cup S_{2}^{\prime}\right)$ is a union of three (non necessarily closed) analytic subsets of dimension $\leqslant s-1$, thus has empty interior in $Z$. It follows that the closure $\bar{Z} \cap\left(B^{\prime} \times B^{\prime \prime}\right)$ is contained in the analytic set $W \subset B^{\prime} \times B^{\prime \prime}$ equal to the common zero set of all functions $P_{\lambda}\left(z^{\prime} ; \lambda\left(z^{\prime \prime}\right)\right)$. Moreover, by construction,

$$
Z \backslash \pi^{-1}\left(A^{\prime} \cup S_{1}^{\prime} \cup S_{2}^{\prime}\right)=W \backslash \pi^{-1}\left(A^{\prime} \cup S_{1}^{\prime} \cup S_{2}^{\prime}\right)
$$

As in the proof of Cor. 5.4, we easily conclude that $\bar{Z} \cap\left(B^{\prime} \times B^{\prime \prime}\right)$ is equal to the union of all irreducible components of $W$ that are not contained in $\pi^{-1}\left(A^{\prime} \cup S_{1}^{\prime} \cup S_{2}^{\prime}\right)$. Hence $\bar{Z}$ is analytic.

Finally, we give two interesting applications of the Remmert-Stein theorem. We assume here that the reader knows what is the complex projective space $\mathbb{P}^{n}$. For more details, see Sect. 5.15.
(8.10) Chow's theorem ([Chow 1949]). Let $A$ be an analytic subset of the complex projective space $\mathbb{P}^{n}$. Then $A$ is algebraic, i.e. $A$ is the common zero set of finitely many homogeneous polynomials $P_{j}\left(z_{0}, \ldots, z_{n}\right), 1 \leqslant j \leqslant N$.

Proof. Let $\pi: \mathbb{C}^{n+1} \backslash\{0\} \longrightarrow \mathbb{P}^{n}$ be the natural projection and $Z=\pi^{-1}(A)$. Then $Z$ is an analytic subset of $\mathbb{C}^{n+1} \backslash\{0\}$ which is invariant by homotheties and $\operatorname{dim} Z=$ $\operatorname{dim} A+1 \geqslant 1$. The Remmert-Stein theorem implies that $\bar{Z}=Z \cup\{0\}$ is an analytic subset of $\mathbb{C}^{n+1}$. Let $f_{1}, \ldots, f_{N}$ be holomorphic functions on a small polydisk $\Delta \subset \mathbb{C}^{n+1}$ of center 0 such that $\bar{Z} \cap \Delta=\bigcap f_{j}^{-1}(0)$. The Taylor series at 0 gives an expansion $f_{j}=\sum_{k=0}^{+\infty} P_{j, k}$ where $P_{j, k}$ is a homogeneous polynomial of degree $k$. We claim that $\bar{Z}$ coincides with the common zero $W$ set of the polynomials $P_{j, k}$. In fact, we clearly have $W \cap \Delta \subset \bigcap f_{j}^{-1}(0)=\bar{Z} \cap \Delta$. Conversely, for $z \in \bar{Z} \cap \Delta$, the invariance of $Z$ by homotheties shows that $f_{j}(t z)=\sum P_{j, k}(z) t^{k}$ vanishes for every complex number $t$ of modulus $<1$, so all coefficients $P_{j, k}(z)$ vanish and $z \in W \cap \Delta$. By homogeneity $\bar{Z}=W$; since $\mathbb{C}\left[z_{0}, \ldots, z_{n}\right]$ is Noetherian, $W$ can be defined by finitely many polynomial equations.
(8.11) E.E. Levi's continuation theorem. Let $X$ be a normal complex space and $A$ an analytic subset such that $\operatorname{dim}(A, x) \leqslant \operatorname{dim}(X, x)-2$ for all $x \in A$. Then every meromorphic function on $X \backslash A$ has a meromorphic extension to $X$.

Proof. We may suppose $X$ irreducible, $\operatorname{dim} X=n$. Let $f$ be a meromorphic function on $X \backslash A$. By Th. 6.13, the pole set $P_{f}$ has pure dimension $(n-1)$, so the Remmert-Stein theorem implies that $\bar{P}_{f}$ is analytic in $X$. Fix a point $x \in A$. There is a connected neighborhood $V$ of $x$ and a non zero holomorphic function $h \in \mathscr{O}_{X}(V)$ such that $\bar{P}_{f} \cap V$ has finitely many irreducible components $\bar{P}_{f, j}$ and $\bar{P}_{f} \cap V \subset h^{-1}(0)$. Select a point $x_{j}$ in $\bar{P}_{f, j} \backslash\left(X_{\text {sing }} \cup\left(\bar{P}_{f}\right)_{\text {sing }} \cup A\right)$. As $x_{j}$ is a regular point on $X$ and on $\bar{P}_{f}$, there is a local coordinate $z_{1, j}$ at $x_{j}$ defining an equation of $\bar{P}_{f, j}$, such that $z_{1, j}^{m_{j}} f \in \mathscr{O}_{X, x_{j}}$ for some integer $m_{j}$. Since $h$ vanishes along $P_{f}$, we have $h^{m_{j}} f \in \mathscr{O}_{X, x}$. Thus, for $m=\max \left\{m_{j}\right\}$, the pole set $P_{g}$ of $g=h^{m} f$ in $V \backslash A$ does not contain $x_{j}$. As $P_{g}$ is $(n-1)$-dimensional and contained in $P_{f} \cap V$, it is a union of irreducible components $\bar{P}_{f, j} \backslash A$. Hence $P_{g}$ must be empty and $g$ is holomorphic on $V \backslash A$. By Cor. 7.8, $g$ has an extension to a holomorphic function $\tilde{g}$ on $V$. Then $\tilde{g} / h^{m}$ is the required meromorphic extension of $f$ on $V$.

## § 9. Complex Analytic Schemes

Our goal is to introduce a generalization of the notion of complex space given in Def. 5.2. A complex space is a space locally isomorphic to an analytic set $A$ in an open subset $\Omega \subset \mathbb{C}^{n}$, together with the sheaf of rings $\mathscr{O}_{A}=\left(\mathscr{O}_{\Omega} / \mathscr{I}_{A}\right)_{\mid A}$. Our desire is to enrich the structure sheaf $\mathscr{O}_{A}$ by replacing $\mathscr{I}_{A}$ with a possibly smaller ideal $\mathscr{f}$ defining the same
zero variety $V(\mathscr{y})=A$. In this way holomorphic functions are described not merely by their values on $A$, but also possibly by some "transversal derivatives" along $A$.

## § 9.1. Ringed Spaces

We start by an abstract notion of ringed space on an arbitrary topological space.
(9.1) Definition. A ringed space is a pair $\left(X, \mathscr{R}_{X}\right)$ consisting of a topological space $X$ and of a sheaf of rings $\mathscr{R}_{X}$ on $X$, called the structure sheaf. A morphism

$$
F:\left(X, \mathscr{R}_{X}\right) \longrightarrow\left(Y, \mathscr{R}_{Y}\right)
$$

of ringed spaces is a pair $\left(f, F^{\star}\right)$ where $f: X \longrightarrow Y$ is a continuous map and

$$
F^{\star}: \quad f^{-1} \mathscr{R}_{Y} \longrightarrow \mathscr{R}_{X}, \quad F_{x}^{\star}: \quad\left(\mathscr{R}_{Y}\right)_{f(x)} \longrightarrow\left(\mathscr{R}_{X}\right)_{x}
$$

a homomorphism of sheaves of rings on $X$, called the comorphism of $F$.

If $F:\left(X, \mathscr{R}_{X}\right) \longrightarrow\left(Y, \mathscr{R}_{Y}\right)$ and $G:\left(Y, \mathscr{R}_{Y}\right) \longrightarrow\left(Z, \mathscr{R}_{Z}\right)$ are morphisms of ringed spaces, the composite $G \circ F$ is the pair consisting of the map $g \circ f: X \longrightarrow Z$ and of the comorphism $(G \circ F)^{\star}=F^{\star} \circ f^{-1} G^{\star}$ :

$$
\begin{align*}
& F^{\star} \circ f^{-1} G^{\star}: f^{-1} g^{-1} \mathscr{R}_{Z} \xrightarrow{f^{-1} G^{\star}} f^{-1} \mathscr{R}_{Y} \quad \xrightarrow{F^{\star}} \mathscr{R}_{X},  \tag{9.2}\\
& F_{x}^{\star} \circ G_{f(x)}^{\star}:\left(\mathscr{R}_{Z}\right)_{g \circ f(x)} \longrightarrow\left(\mathscr{R}_{X}\right)_{x} .
\end{align*}
$$

## § 9.2. Definition of Complex Analytic Schemes

We begin by a description of what will be the local model of an analytic scheme. Let $\Omega \subset \mathbb{C}^{n}$ be an open subset, $\mathscr{G} \subset \mathbb{O}_{\Omega}$ a coherent sheaf of ideals and $A=V(\mathscr{F})$ the analytic set in $\Omega$ defined locally as the zero set of a system of generators of $\mathscr{g}$. By Hilbert's Nullstellensatz 4.22 we have $\mathscr{I}_{A}=\sqrt{\mathscr{g}}$, but $\mathscr{I}_{A}$ differs in general from $\mathscr{g}$. The sheaf of rings $\mathbb{O}_{\Omega} / \mathscr{G}$ is supported on $A$, i.e. $\left(\mathscr{O}_{\Omega} / \mathscr{F}\right)_{x}=0$ if $x \notin A$. Ringed spaces of the type $\left(A, \oslash_{\Omega} / \mathscr{F}\right)$ will be used as the local models of analytic schemes.
(9.3) Definition. A morphism

$$
F=\left(f, F^{\star}\right):\left(A, \mathscr{O}_{\Omega} / \mathscr{F}_{卜 A}\right) \longrightarrow\left(A^{\prime}, \mathscr{O}_{\Omega^{\prime}} / \mathscr{F}_{\uparrow A^{\prime}}^{\prime}\right)
$$

is said to be analytic if for every point $x \in A$ there exists a neighborhood $W_{x}$ of $x$ in $\Omega$ and a holomorphic function $\Phi: W_{x} \longrightarrow \Omega^{\prime}$ such that $f_{\uparrow A \cap W_{x}}=\Phi_{\left\lceil A \cap W_{x}\right.}$ and such that the comorphism

$$
F_{x}^{\star}:\left(\mathfrak{O}_{\Omega^{\prime}} / \mathscr{g}^{\prime}\right)_{f(x)} \longrightarrow\left(\mathbb{O}_{\Omega} / \mathscr{F}\right)_{x}
$$

is induced by $\Phi^{\star}: \mathbb{O}_{\Omega^{\prime}, f(x)} \ni u \longmapsto u \circ \Phi \in \mathbb{O}_{\Omega, x}$ with $\Phi^{\star} \mathscr{g}^{\prime} \subset \mathscr{g}$.
(9.4) Example. Take $\Omega=\mathbb{C}^{n}$ and $\mathscr{g}=\left(z_{n}^{2}\right)$. Then $A$ is the hyperplane $\mathbb{C}^{n-1} \times\{0\}$, and the sheaf $\mathbb{O}_{\mathbb{C}^{n}} / \mathscr{g}$ can be identified with the sheaf of rings of functions $u+z_{n} u^{\prime}$, $u, u^{\prime} \in \mathscr{O}_{\mathbb{C}^{n-1}}$, with the relation $z_{n}^{2}=0$. In particular, $z_{n}$ is a nilpotent element of $\mathscr{O}_{\mathbb{C}^{n}} / \mathscr{g}$. A morphism $F$ of $\left(A, \mathscr{O}_{\mathbb{C}^{n}} / \mathscr{F}\right)$ into itself is induced (at least locally) by a holomorphic
map $\Phi=\left(\widetilde{\Phi}, \Phi_{n}\right)$ defined on a neighborhood of $A$ in $\mathbb{C}^{n}$ with values in $\mathbb{C}^{n}$, such that $\Phi(A) \subset A$, i.e. $\Phi_{n \upharpoonright A}=0$. We see that $F$ is completely determined by the data

$$
\begin{aligned}
f\left(z_{1}, \ldots, z_{n-1}\right) & =\widetilde{\Phi}\left(z_{1}, \ldots, z_{n-1}, 0\right), & f: & \mathbb{C}^{n-1} \longrightarrow \mathbb{C}^{n-1} \\
f^{\prime}\left(z_{1}, \ldots, z_{n-1}\right) & =\frac{\partial \Phi}{\partial z_{n}}\left(z_{1}, \ldots, z_{n-1}, 0\right), & f^{\prime}: & \mathbb{C}^{n-1} \longrightarrow \mathbb{C}^{n}
\end{aligned}
$$

which can be chosen arbitrarily.
(9.5) Definition. A complex analytic scheme is a ringed space ( $X, \mathscr{O}_{X}$ ) over a separable Hausdorff topological space $X$, satisfying the following property: there exist an open covering $\left(U_{\lambda}\right)$ of $X$ and isomorphisms of ringed spaces

$$
G_{\lambda}:\left(U_{\lambda}, \mathscr{O}_{X \mid U_{\lambda}}\right) \longrightarrow\left(A_{\lambda}, \mathscr{O}_{\Omega_{\lambda}} / \mathscr{H}_{\lambda \upharpoonright A_{\lambda}}\right)
$$

where $A_{\lambda}$ is the zero set of a coherent sheaf of ideals $\mathscr{F}_{\lambda}$ on an open subset $\Omega_{\lambda} \subset$ $\mathbb{C}^{N_{\lambda}}$, such that every transition morphism $G_{\lambda} \circ G_{\mu}^{-1}$ is a holomorphic isomorphism from $g_{\mu}\left(U_{\lambda} \cap U_{\mu}\right) \subset A_{\mu}$ onto $g_{\lambda}\left(U_{\lambda} \cap U_{\mu}\right) \subset A_{\lambda}$, equipped with the respective structure sheaves $\mathfrak{O}_{\Omega_{\mu}} / \mathscr{F}_{\mu \upharpoonright A_{\mu}}, \mathfrak{O}_{\Omega_{\lambda}} / \mathscr{F}_{\lambda \upharpoonright A_{\lambda}}$.

We shall often consider the maps $G_{\lambda}$ as identifications and write simply $U_{\lambda}=A_{\lambda}$. A morphism $F:\left(X, \mathscr{O}_{X}\right) \longrightarrow\left(Y, \mathscr{O}_{Y}\right)$ of analytic schemes obtained by gluing patches $\left(A_{\lambda}, \mathscr{O}_{\Omega_{\lambda}} / \mathscr{F}_{\lambda \mid A_{\lambda}}\right)$ and $\left(A_{\mu}^{\prime}, \mathscr{O}_{\Omega_{\mu}^{\prime}} / \mathscr{I}_{\mu A_{\mu}^{\prime}}^{\prime}\right)$, respectively, is a morphism $F$ of ringed spaces such that for each pair $(\lambda, \mu)$, the restriction of $F$ from $A_{\lambda} \cap f^{-1}\left(A_{\mu}^{\prime}\right) \subset X$ to $A_{\mu}^{\prime} \subset Y$ is holomorphic in the sense of Def. 9.3.

## § 9.3. Nilpotent Elements and Reduced Schemes

Let $\left(X, \mathscr{O}_{X}\right)$ be an analytic scheme. The set of nilpotent elements is the sheaf of ideals of $\mathscr{O}_{X}$ defined by

$$
\begin{equation*}
\mathcal{N}_{X}=\left\{u \in \mathbb{O}_{X} ; u^{k}=0 \text { for some } k \in \mathbb{N}\right\} \tag{9.6}
\end{equation*}
$$

Locally, we have $\mathscr{O}_{X \upharpoonright A_{\lambda}}=\left(\mathscr{O}_{\Omega_{\lambda}} / \mathscr{I}_{\lambda}\right)_{\mid A_{\lambda}}$, thus

$$
\begin{align*}
& \mathcal{N}_{X \mid A_{\lambda}}=\left(\sqrt{\mathscr{g}_{\lambda}} / \mathscr{g}_{\lambda}\right)_{\mid A_{\lambda}}  \tag{9.7}\\
& \left(\mathscr{O}_{X} / \mathcal{N}_{X}\right)_{\mid A_{\lambda}} \simeq\left(\mathscr{O}_{\Omega_{\lambda}} / \sqrt{\mathscr{H}_{\lambda}}\right)_{\mid A_{\lambda}}=\left(\mathscr{O}_{\Omega_{\lambda}} / \mathscr{I}_{A_{\lambda}}\right)_{\mid A_{\lambda}}=\mathscr{O}_{A_{\lambda}} \tag{9.8}
\end{align*}
$$

The scheme $\left(X, \mathscr{O}_{X}\right)$ is said to be reduced if $\mathcal{N}_{X}=0$. The associated ringed space $\left(X, \mathscr{O}_{X} / \mathcal{N}_{X}\right)$ is reduced by construction; it is called the reduced scheme of $\left(X, \mathscr{O}_{X}\right)$. We shall often denote the original scheme by the letter $X$ merely, the associated reduced scheme by $X_{\text {red }}$, and let $\mathscr{O}_{X, \text { red }}=\mathscr{O}_{X} / \mathcal{N}_{X}$. There is a canonical morphism $X_{\text {red }} \rightarrow X$ whose comorphism is the reduction morphism

$$
\begin{equation*}
\mathfrak{O}_{X}(U) \longrightarrow \mathfrak{O}_{X, \text { red }}(U)=\left(\mathscr{O}_{X} / \mathcal{N}_{X}\right)(U), \quad \forall U \text { open set in } X \tag{9.9}
\end{equation*}
$$

By (9.8), the notion of reduced scheme is equivalent to the notion of complex space introduced in Def. 5.2. It is easy to see that a morphism $F$ of reduced schemes $X, Y$ is completely determined by the set-theoretic map $f: X \longrightarrow Y$.

## § 9.4. Coherent Sheaves on Analytic Schemes

If $\left(X, \mathscr{O}_{X}\right)$ is an analytic scheme, a sheaf $\mathscr{S}$ of $\mathscr{O}_{X}$-modules is said to be coherent if it satisfies the same properties as those already stated when $X$ is a manifold:
(9.10) $\mathscr{S}$ is locally finitely generated over $\mathscr{O}_{X}$;
(9.10') for any open set $U \subset X$ and any sections $G_{1}, \ldots, G_{q} \in \mathscr{P}(U)$, the relation sheaf $\mathscr{R}\left(G_{1}, \ldots, G_{q}\right) \subset \mathscr{O}_{X \mid U}^{\oplus q}$ is locally finitely generated.
Locally, we have $\mathscr{O}_{X \upharpoonright A_{\lambda}}=\mathscr{O}_{\Omega_{\lambda}} / \mathscr{F}_{\lambda}$, so if $i_{\lambda}: A_{\lambda} \rightarrow \Omega_{\lambda}$ is the injection, the direct image $\mathscr{S}_{\lambda}=\left(i_{\lambda}\right)_{\star}\left(\mathscr{S}_{\uparrow A_{\lambda}}\right)$ is a module over $\mathscr{O}_{\Omega_{\lambda}}$ such that $\mathscr{g}_{\lambda} \cdot \mathscr{S}_{\lambda}=0$. It is clear that $\mathscr{S}_{\mid \Omega_{\lambda}}$ is coherent if and only if $\mathscr{S}_{\lambda}$ is coherent as a module over $\mathscr{Q}_{\Omega_{\lambda}}$. It follows immediately that the Oka theorem and its consequences 3.16-20 are still valid over analytic schemes.

## $\S$ 9.5. Subschemes

Let $X$ be an analytic scheme and $\mathscr{G}$ a coherent sheaf of ideals in $\mathscr{Q}_{X}$. The image of $\mathscr{G}_{\mathcal{G}}$ in $\mathscr{O}_{X, \text { red }}$ is a coherent sheaf of ideals, and its zero set $Y$ is clearly an analytic subset of $X_{\text {red }}$. We can make $Y$ into a scheme by introducing the structure sheaf

$$
\begin{equation*}
\mathfrak{O}_{Y}=\left(\mathscr{O}_{X} / \mathscr{G}\right)_{\mid Y} \tag{9.11}
\end{equation*}
$$

and we have a scheme morphism $F:\left(Y, \mathscr{O}_{Y}\right) \longrightarrow\left(X, \mathscr{O}_{X}\right)$ such that $f$ is the inclusion and $F^{\star}: f^{-1} \mathscr{O}_{X} \longrightarrow \mathscr{O}_{Y}$ the obvious map of $\mathscr{O}_{X \upharpoonright Y}$ onto its quotient $\mathscr{O}_{Y}$. The scheme $\left(Y, \mathscr{O}_{Y}\right)$ will be denoted $V(\mathscr{G})$. When the analytic set $Y$ is given, the structure sheaf of $V(\mathscr{G})$ depends of course on the choice of the equations of $Y$ in the ideal $\mathscr{G}$; in general $\mathscr{O}_{Y}$ has nilpotent elements.

## $\S$ 9.6. Inverse Images of Coherent Sheaves

Let $F:\left(X, \mathscr{®}_{X}\right) \longrightarrow\left(Y, \mathscr{O}_{Y}\right)$ be a morphism of analytic schemes and $\mathscr{S}$ a coherent sheaf over $Y$. The sheaf theoretic inverse image $f^{-1} \mathscr{S}$, whose stalks are $\left(f^{-1} \mathscr{S}\right)_{x}=\mathscr{S}_{f(x)}$, is a sheaf of modules over $f^{-1} \mathscr{O}_{Y}$. We define the analytic inverse image $F^{\star} \mathscr{S}$ by

$$
\begin{equation*}
F^{\star} \mathscr{S}=\mathscr{O}_{X} \otimes_{f^{-1} \mathscr{O}_{Y}} f^{-1} \mathscr{S}, \quad\left(F^{\star} \mathscr{S}\right)_{x}=\mathscr{O}_{X, x} \otimes_{\mathscr{Q}_{Y, f(x)}} \mathscr{T}_{f(x)} \tag{9.12}
\end{equation*}
$$

Here the tensor product is taken with respect to the comorphism $F^{\star}: f^{-1} \mathscr{O}_{Y} \rightarrow \mathscr{O}_{X}$, which yields a ring morphism $\mathscr{O}_{Y, f(x)} \rightarrow \mathscr{O}_{X, x}$. If $\mathscr{S}$ is given over $U \subset Y$ by a local presentation

$$
\mathfrak{O}_{Y \upharpoonright U}^{\oplus p} \xrightarrow{A} \mathfrak{O}_{Y \upharpoonright U}^{\oplus q} \longrightarrow \mathscr{S}_{\uparrow U} \longrightarrow 0
$$

where $A$ is a $(q \times p)$-matrix with coefficients in $\mathscr{C}_{Y}(U)$, our definition shows that $F^{\star} \mathscr{S}$ is a coherent sheaf over $\mathscr{O}_{X}$, given over $f^{-1}(U)$ by the local presentation

$$
\begin{equation*}
\mathscr{O}_{X \mid f^{-1}(U)}^{\oplus p} \xrightarrow{F^{\star} A} \mathscr{O}_{X \mid f^{-1}(U)}^{\oplus q} \longrightarrow F^{\star} \mathscr{S}_{\mid f-1}(U) \longrightarrow 0 . \tag{9.13}
\end{equation*}
$$

## § 9.7. Products of Analytic Schemes

Let $\left(X, \mathscr{O}_{X}\right)$ and $\left(Y, \mathscr{O}_{Y}\right)$ be analytic schemes, and let $\left(A_{\lambda}, \mathscr{O}_{\Omega_{\lambda}} / \mathscr{F}_{\lambda}\right),\left(B_{\mu}, \mathscr{O}_{\Omega_{\mu}^{\prime}} / \mathscr{F}_{\mu}^{\prime}\right)$ be local models of $X, Y$, respectively. The product scheme $\left(X \times Y, 0_{X \times Y}\right)$ is obtained by gluing the open patches

$$
\begin{equation*}
\left(A_{\lambda} \times B_{\mu}, \mathscr{O}_{\Omega_{\lambda} \times \Omega_{\mu}^{\prime}} /\left(\operatorname{pr}_{1}^{-1} \mathscr{G}_{\lambda}+\operatorname{pr}_{2}^{-1} \mathscr{G}_{\mu}^{\prime}\right) \mathscr{O}_{\Omega_{\lambda} \times \Omega_{\mu}^{\prime}}\right) . \tag{9.14}
\end{equation*}
$$

In other words, if $A_{\lambda}, B_{\mu}$ are the subschemes of $\Omega_{\lambda}, \Omega_{\mu}^{\prime}$ defined by the equations $g_{\lambda, j}(x)=$ $0, g_{\mu, k}^{\prime}(y)=0$, where $\left(g_{\lambda, j}\right)$ and $\left(g_{\mu, k}^{\prime}\right)$ are generators of $g_{\lambda}$ and $\mathscr{g}_{\mu}^{\prime}$ respectively, then $A_{\lambda} \times B_{\mu}$ is equipped with the structure sheaf $\mathbb{O}_{\Omega_{\lambda} \times \Omega_{\mu}^{\prime}} /\left(g_{\lambda, j}(x), g_{\mu, k}^{\prime}(y)\right)$.

Now, let $\mathscr{S}$ be a coherent sheaf over $\mathscr{O}_{X}$ and let $\mathscr{G}^{\prime}$ be a coherent sheaf over $\mathscr{O}_{Y}$. The (analytic) external tensor product $\mathscr{S} \boxtimes \mathscr{S}^{\prime}$ is defined to be

$$
\begin{equation*}
\mathscr{S} \boxtimes \mathscr{S}^{\prime}=\operatorname{pr}_{1}^{\star} \mathscr{S} \otimes_{\mathbb{Q}_{X \times Y}} \operatorname{pr}_{2}^{\star} \mathscr{S}^{\prime} \tag{9.15}
\end{equation*}
$$

If we go back to the definition of the inverse image, we see that the stalks of $\mathscr{S}_{\boxtimes} \mathscr{S}^{\prime}$ are given by

$$
\left(\mathscr{S}^{\otimes} \mathscr{S}^{\prime}\right)_{(x, y)}=\mathscr{O}_{X \times Y,(x, y)} \otimes_{\mathbb{Q}_{X, x} \otimes \mathbb{Q}_{Y, y}}\left(\mathscr{S}_{x} \otimes_{\mathbb{C}} \mathscr{S}_{y}^{\prime}\right)
$$

in particular $\left(\mathscr{S} \boxtimes \mathscr{S}^{\prime}\right)_{(x, y)}$ does not coincide with the sheaf theoretic tensor product $\mathscr{S}_{x} \otimes$ $\mathscr{S}_{y}^{\prime}$ which is merely a module over $\mathscr{O}_{X, x} \otimes \mathscr{O}_{Y, y}$. If $\mathscr{S}$ and $\mathscr{S}^{\prime}$ are given by local presentations

$$
\mathscr{O}_{X \upharpoonright U}^{\oplus p} \xrightarrow{A} \mathscr{O}_{X \upharpoonright U}^{\oplus q} \longrightarrow \mathscr{S}_{\uparrow U} \longrightarrow 0, \quad \mathscr{O}_{Y \upharpoonright U^{\prime}}^{p^{\prime}} \xrightarrow{B} \mathscr{O}_{Y \mid U^{\prime}}^{q^{\prime}} \longrightarrow \mathscr{S}_{\upharpoonright U^{\prime}}^{\prime} \longrightarrow 0,
$$

then $\mathscr{S} \boxtimes \mathscr{I}^{\prime}$ is the coherent sheaf given by

$$
\mathscr{O}_{X \times Y \mid U \times U^{\prime}}^{p{ }^{\prime} \oplus q p^{\prime}} \xrightarrow{(A(x) \otimes \mathrm{Id}, \mathrm{Id} \otimes B(y))} \mathscr{O}_{X \times Y \upharpoonright U \times U^{\prime}}^{q q^{\prime}} \longrightarrow\left(\mathscr{S}^{\prime} \boxtimes \mathscr{S}^{\prime}\right)_{\mid U \times U^{\prime}} \longrightarrow 0 .
$$

## § 9.8. Zariski Embedding Dimension

If $x$ is a point of an analytic scheme $\left(X, \mathscr{O}_{X}\right)$, the Zariski embedding dimension of the germ $(X, x)$ is the smallest integer $N$ such that $(X, x)$ can be embedded in $\mathbb{C}^{N}$, i.e. such that there exists a patch of $X$ near $x$ isomorphic to $\left(A, \sigma_{\Omega} / \mathscr{J}\right)$ where $\Omega$ is an open subset of $\mathbb{C}^{N}$. This dimension is denoted

$$
\begin{equation*}
\operatorname{embdim}(X, x)=\text { smallest such } N . \tag{9.16}
\end{equation*}
$$

Consider the maximal ideal $\mathfrak{m}_{X, x} \subset \mathscr{O}_{X, x}$ of functions which vanish at point $x$. If ( $X, x$ ) is embedded in $(\Omega, x)=\left(\mathbb{C}^{N}, 0\right)$, then $\mathfrak{m}_{X, x} / \mathfrak{m}_{X, x}^{2}$ is generated by $z_{1}, \ldots, z_{N}$, so $d=$ $\operatorname{dim} \mathfrak{m}_{X, x} / \mathfrak{m}_{X, x}^{2} \leqslant N$. Let $s_{1}, \ldots, s_{d}$ be germs in $\mathfrak{m}_{\Omega, x}$ which yield a basis of $\mathfrak{m}_{X, x} / \mathfrak{m}_{X, x}^{2} \simeq$ $\mathfrak{m}_{\Omega, x} /\left(\mathfrak{m}_{\Omega, x}^{2}+\mathscr{I}_{x}\right)$. We can write

$$
z_{j}=\sum_{1 \leqslant k \leqslant d} c_{j k} s_{k}+u_{j}+f_{j}, \quad c_{j k} \in \mathbb{C}, \quad u_{j} \in \mathfrak{m}_{\Omega, x}^{2}, \quad f_{j} \in \mathscr{g}_{x}, \quad 1 \leqslant j \leqslant n
$$

Then we find $d z_{j}=\sum c_{j k} d s_{k}(x)+d f_{j}(x)$, so that the rank of the system of differentials $\left(d f_{j}(x)\right)$ is at least $N-d$. Assume for example that $d f_{1}(x), \ldots, d f_{N-d}(x)$ are linearly independant. By the implicit function theorem, the equations $f_{1}=\ldots=f_{N-d}=0$ define a germ of smooth subvariety $(Z, x) \subset(\Omega, x)$ of dimension $d$ which contains $(X, x)$. We have $\mathscr{O}_{Z}=\mathscr{O}_{\Omega} /\left(f_{1}, \ldots, f_{N-d}\right)$ in a neighborhood of $x$, thus

$$
\mathfrak{O}_{X}=\mathscr{O}_{\Omega} / \mathscr{G} \simeq \mathscr{O}_{Z} / \mathscr{g}^{\prime} \quad \text { where } \mathscr{g}^{\prime}=\mathscr{g} /\left(f_{1}, \ldots, f_{N-d}\right)
$$

This shows that $(X, x)$ can be imbedded in $\mathbb{C}^{d}$, and we get

$$
\begin{equation*}
\operatorname{embdim}(X, x)=\operatorname{dim} \mathfrak{m}_{X, x} / \mathfrak{m}_{X, x}^{2} \tag{9.17}
\end{equation*}
$$

(9.18) Remark. For a given dimension $n=\operatorname{dim}(X, x)$, the embedding dimension $d$ can be arbitrarily large. Consider for example the curve $\Gamma \subset \mathbb{C}^{N}$ parametrized by $\mathbb{C} \ni t \longmapsto\left(t^{N}, t^{N+1}, \ldots, t^{2 N-1}\right)$. Then $\mathbb{O}_{\Gamma, 0}$ is the ring of convergent series in $\mathbb{C}\{t\}$ which have no terms $t, t^{2}, \ldots, t^{N-1}$, and $\mathfrak{m}_{\Gamma, 0} / \mathfrak{m}_{\Gamma, 0}^{2}$ admits precisely $z_{1}=t^{N}, \ldots, z_{N}=t^{2 N-1}$ as a basis. Therefore $n=1$ but $d=N$ can be as large as we want.

## $\S$ 10. Bimeromorphic maps, Modifications and Blow-ups

It is a very frequent situation in analytic or algebraic geometry that two complex spaces have isomorphic dense open subsets but are nevertheless different along some analytic subset. These ideas are made precise by the notions of modification and bimeromorphic map. This will also lead us to generalize the notion of meromorphic function to maps between analytic schemes. If $\left(X, \mathscr{O}_{X}\right)$ is an analytic scheme, $\mathscr{M}_{X}$ denotes the sheaf of meromorphic functions on $X$, defined at the beginning of § 6.2 .
(10.1) Definition. Let $\left(X, \mathscr{O}_{X}\right),\left(Y, \mathscr{O}_{Y}\right)$ be analytic schemes. An analytic morphism $F: X \rightarrow Y$ is said to be a modification if $F$ is proper and if there exists a nowhere dense closed analytic subset $B \subset Y$ such that the restriction $F: X \backslash F^{-1}(B) \rightarrow Y \backslash B$ is an isomorphism.
(10.2) Definition. If $F: X \rightarrow Y$ is a modification, then the comorphism $F^{\star}: f^{\star} \mathscr{O}_{Y} \rightarrow$ $\mathfrak{O}_{X}$ induces an isomorphism $F^{\star}: f^{\star} \mathcal{M}_{Y} \rightarrow \mathcal{M}_{X}$ for the sheaves of meromorphic functions on $X$ and $Y$.

Proof. Let $v=g / h$ be a section of $\mathcal{M}_{Y}$ on a small open set $\Omega$ where $u$ is actually given as a quotient of functions $g, h \in \mathscr{C}_{Y}(\Omega)$. Then $F^{\star} u=(g \circ F) /(h \circ F)$ is a section of $\mathcal{M}_{X}$ on $F^{-1}(\Omega)$, for $h \circ F$ cannot vanish identically on any open subset $W$ of $F^{-1}(\Omega)$ (otherwise $h$ would vanish on the open subset $F\left(W \backslash F^{-1}(B)\right)$ of $\left.\Omega \backslash B\right)$. Thus the extension of the comorphism to sheaves of meromorphic functions is well defined. Our claim is that this is an isomorphism. The injectivity of $F^{\star}$ is clear: $F^{\star} u=0$ implies $g \circ F=0$, which implies $g=0$ on $\Omega \backslash B$ and thus $g=0$ on $\Omega$ because $B$ is nowhere dense. In order to prove surjectivity, we need only show that every section $u \in \mathscr{O}_{X}\left(F^{-1}(\Omega)\right)$ is in the image of $\mathcal{M}_{Y}(\Omega)$ by $F^{\star}$. For this, we may shrink $\Omega$ into a relatively compact subset $\Omega^{\prime} \subset \subset \Omega$ and thus assume that $u$ is bounded (here we use the properness of $F$ through the fact that $F^{-1}\left(\Omega^{\prime}\right)$ is relatively compact in $\left.F^{-1}(\Omega)\right)$. Then $v=u \circ F^{-1}$ defines a bounded holomorphic function on $\Omega \backslash B$. By Th. 7.2, it follows that $v$ is weakly holomorphic for the reduced structure of $Y$. Our claim now follows from the following Lemma.
(10.3) Lemma. If $\left(X, \mathscr{O}_{X}\right)$ is an analytic scheme, then every holomorphic function $v$ in the complement of a nowhere dense analytic subset $B \subset Y$ which is weakly holomorphic on $X_{\text {red }}$ is meromorphic on $X$.

Proof. It is enough to argue with the germ of $v$ at any point $x \in Y$, and thus we may suppose that $\left(Y, \mathscr{O}_{Y}\right)=\left(A, \mathscr{O}_{\Omega} / \mathscr{I}\right)$ is embedded in $\mathbb{C}^{N}$. Because $v$ is weakly holomorphic, we can write $v=g / h$ in $Y_{\text {red }}$, for some germs of holomorphic functions $g, h$. Let $\widetilde{g}$ and $\widetilde{h}$ be extensions of $g, h$ to $\mathbb{O}_{\Omega, x}$. Then there is a neighborhood $U$ of $x$ such that $\widetilde{g}-v \widetilde{h}$ is a nilpotent section of $c O_{\Omega}(U \backslash B)$ which is in $\mathscr{I}$ on
(10.4) Definition. A meromorphic map $F: X-\rightarrow Y$ is a scheme morphism $F$ : $X \backslash A \rightarrow Y$ defined in the complement of a nowhere dense analytic subset $A \subset X$, such that the closure of the graph of $F$ in $X \times Y$ is an analytic subset (for the reduced complex space structure of $X \times Y$ ).

## § 11. Exercises

$\oint$ 11.1. Let $\mathscr{A}$ be a sheaf on a topological space $X$. If the sheaf space $\widetilde{\mathscr{A}}$ is Hausdorff, show that $\mathscr{A}$ satisfies the following unique continuation principle: any two sections $s, s^{\prime} \in \mathscr{A}(U)$ on a connected open set $U$ which coincide on some non empty open subset $V \subset U$ must coincide identically on $U$. Show that the converse holds if $X$ is Hausdorff and locally connected.
$\oint$ 11.2. Let $\mathscr{A}$ be a sheaf of abelian groups on $X$ and let $s \in \mathscr{A}(X)$. The support of $s$, denoted Supp $s$, is defined to be $\{x \in X ; s(x) \neq 0\}$. Show that $\operatorname{Supp} s$ is a closed subset of $X$. The support of $\mathscr{A}$ is defined to be $\operatorname{Supp} \mathscr{A}=\left\{x \in X ; \mathcal{A}_{x} \neq 0\right\}$. Show that Supp $\mathscr{A}$ is not necessarily closed: if $\Omega$ is an open set in $X$, consider the sheaf $\mathscr{A}$ such that $\mathscr{A}(U)$ is the set of continuous functions $f \in \mathscr{C}(U)$ which vanish on a neighborhood of $U \cap(X \backslash \Omega)$.
$\oint$ 11.3. Let $\mathscr{A}$ be a sheaf of rings on a topological space $X$ and let $\mathscr{F}, \mathscr{G}$ be sheaves of $\mathscr{A}$-modules. We define a presheaf $\mathscr{H}=\mathscr{H} \operatorname{om}_{\mathscr{A}}(\mathscr{F}, \mathscr{G})$ such that $\mathscr{H}(U)$ is the module of all sheaf-homomorphisms $\mathscr{F}_{\uparrow U} \rightarrow \mathscr{G}_{\uparrow U}$ which are $\mathscr{A}$-linear.
a) Show that $\mathscr{H} \operatorname{om}_{\mathscr{A}}(\mathscr{F}, \mathscr{G})$ is a sheaf and that there exists a canonical homomorphism $\varphi_{x}: \mathscr{H} \operatorname{om}_{\mathscr{A}}(\mathscr{F}, \mathscr{G})_{x} \longrightarrow \operatorname{hom}_{\mathfrak{d}_{x}}\left(\mathscr{F} x, \varphi_{x}\right)$ for every $x \in X$.
b) If $\mathscr{F}$ is locally finitely generated, then $\varphi_{x}$ is injective, and if $\mathscr{F}$ has local finite presentations as in (3.12), then $\varphi_{x}$ is bijective.
c) Suppose that $\mathscr{A}$ is a coherent sheaf of rings and that $\mathscr{F}, \mathscr{G}$ are coherent modules over $\mathscr{A}$. Then $\mathscr{H} \operatorname{om}_{\mathscr{A}}\left(\mathscr{F}, \mathscr{G}_{)}\right)$is a coherent $\mathscr{A}$-module.
Hint: observe that the result is true if $\mathscr{F}=\mathscr{A}^{\oplus p}$ and use a local presentation of $\mathscr{F}$ to get the conclusion.
$\oint$ 11.4. Let $f: X \rightarrow Y$ be a continuous map of topological spaces. Given sheaves of abelian groups $\mathscr{A}$ on $X$ and $\mathscr{B}$ on $Y$, show that there is a natural isomorphism

$$
\operatorname{hom}_{X}\left(f^{-1} \mathscr{B}, \mathscr{A}\right)=\operatorname{hom}_{Y}\left(\mathscr{B}, f_{\star} \mathscr{A}\right)
$$

Hint: use the natural morphisms (2.17).
$\oint$ 11.5. Show that the sheaf of polynomials over $\mathbb{C}^{n}$ is a coherent sheaf of rings (with either the ordinary topology or the Zariski topology on $\mathbb{C}^{n}$ ). Extend this result to the case of regular algebraic functions on an algebraic variety.
Hint: check that the proof of the Oka theorem still applies.
$\oint$ 11.6. Let $P$ be a non zero polynomial on $\mathbb{C}^{n}$. If $P$ is irreducible in $\mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$, show that the hypersurface $H=P^{-1}(0)$ is globally irreducible as an analytic set. In general, show that the irreducible components of $H$ are in a one-to-one correspondence with the irreducible factors of $P$.
Hint: for the first part, take coordinates such that $P\left(0, \ldots, 0, z_{n}\right)$ has degree equal to $P$; if $H$ splits in two components $H_{1}, H_{2}$, then $P$ can be written as a product $P_{1} P_{2}$ where the roots of $P_{j}\left(z^{\prime}, z_{n}\right)$ correspond to points in $H_{j}$.
$\S$ 11.7. Prove the following facts:
a) For every algebraic variety $A$ of pure dimension $p$ in $\mathbb{C}^{n}$, there are coordinates $z^{\prime}=\left(z_{1}, \ldots, z_{p}\right)$, $z^{\prime \prime}=\left(z_{p+1}, \ldots, z_{n}\right)$ such that $\pi: A \rightarrow \mathbb{C}^{p}, z \mapsto z^{\prime \prime}$ is proper with finite fibers, and such that $A$ is entirely contained in a cone

$$
\left|z^{\prime \prime}\right| \leqslant C\left(\left|z^{\prime}\right|+1\right)
$$

Hint: imitate the proof of Cor. 4.11.
b) Conversely if an analytic set $A$ of pure dimension $p$ in $\mathbb{C}^{n}$ is contained in a cone $\left|z^{\prime \prime}\right| \leqslant C\left(\left|z^{\prime}\right|+1\right)$, then $A$ is algebraic.
Hint: first apply (5.9) to conclude that the projection $\pi: A \rightarrow \mathbb{C}^{p}$ is finite. Then repeat the arguments used in the final part of the proof of Th. 4.19.
c) Deduce from a), b) that an algebraic set in $\mathbb{C}^{n}$ is irreducible if and only if it is irreducible as an analytic set.
$\oint$ 11.8. Let $\Gamma: f(x, y)=0$ be a germ of analytic curve in $\mathbb{C}^{2}$ through $(0,0)$ and let $\left(\Gamma_{j}, 0\right)$ be the irreducible components of $(\Gamma, 0)$. Suppose that $f(0, y) \not \equiv 0$. Show that the roots $y$ of $f(x, y)=0$ corresponding to points of $\Gamma$ near 0 are given by Puiseux expansions of the form $y=g_{j}\left(x^{1 / q_{j}}\right)$, where $g_{j} \in \mathscr{O}_{\mathbb{C}, 0}$ and where $q_{j}$ is the sheet number of the projection $\Gamma_{j} \rightarrow \mathbb{C},(x, y) \mapsto x$.
Hint: special case of the parametrization obtained in (4.27).
$\oint$ 11.9. The goal of this exercise is to prove the existence and the analyticity of the tangent cone to an arbitrary analytic germ $(A, 0)$ in $\mathbb{C}^{n}$. Suppose that $A$ is defined by holomorphic equations $f_{1}=\ldots=$ $f_{N}=0$ in a ball $\Omega=B(0, r)$. Then the (set theoretic) tangent cone to $A$ at 0 is the set $C(A, 0)$ of all limits of sequences $t_{\nu}^{-1} z_{\nu}$ with $z_{\nu} \in A$ and $\mathbb{C}^{\star} \ni t_{\nu}$ converging to 0 .
a) Let $E$ be the set of points $(z, t) \in \Omega \times \mathbb{C}^{\star}$ such that $z \in t^{-1} A$. Show that the closure $\bar{E}$ in $\Omega \times \mathbb{C}$ is analytic.
Hint: observe that $E=A \backslash(\Omega \times\{0\})$ where $A=\left\{f_{j}(t z)=0\right\}$ and apply Cor. 5.4.
b) Show that $C(A, 0)$ is a conic set and that $\bar{E} \cap(\Omega \times\{0\})=C(A, 0) \times\{0\}$ and conclude. Infer from this that $C(A, 0)$ is an algebraic subset of $\mathbb{C}^{n}$.
$\oint$ 11.10. Give a new proof of Theorem 5.5 based on the coherence of ideal sheaves and on the strong noetherian property.
$\oint$ 11.11. Let $X$ be an analytic space and let $A, B$ be analytic subsets of pure dimensions. Show that $\operatorname{codim}_{X}(A \cap B) \leqslant \operatorname{codim}_{X} A+\operatorname{codim}_{X} B$ if $A$ or $B$ is a local complete intersection, but that the equality does not necessarily hold in general.
Hint: see Remark (6.5).
$\oint$ 11.12. Let $\Gamma$ be the curve in $\mathbb{C}^{3}$ parametrized by $\mathbb{C} \ni t \longmapsto(x, y, z)=\left(t^{3}, t^{4}, t^{5}\right)$. Show that the ideal sheaf $\mathscr{I}_{\Gamma}$ is generated by the polynomials $x z-y^{2}, x^{3}-y z$ and $x^{2} y-z^{2}$, and that the germ $(\Gamma, 0)$ is not a (sheaf theoretic) complete intersection.
Hint: $\Gamma$ is smooth except at the origin. Let $f(x, y, z)=\sum a_{\alpha \beta \gamma} x^{\alpha} y^{\beta} z^{\gamma}$ be a convergent power series near 0 . Show that $f \in \mathscr{I}_{\Gamma, 0}$ if and only if all weighted homogeneous components

$$
f_{k}=\sum_{3 \alpha+4 \beta+5 \gamma=k} a_{\alpha \beta \gamma} x^{\alpha} y^{\beta} z^{\gamma}
$$

are in $\mathscr{I}_{\Gamma, 0}$. By means of suitable substitutions, reduce the proof to the case when $f=f_{k}$ is homogeneous with all non zero monomials satisfying $\alpha \leqslant 2, \beta \leqslant 1, \gamma \leqslant 1$; then check that there is at most one such monomial in each weighted degree $\leqslant 15$ the product of a power of $x$ by a homogeneous polynomial of weighted degree $\leqslant 8$.

## Chapter III

## Positive Currents and Lelong Numbers

In 1957, P. Lelong introduced natural positivity concepts for currents of pure bidimension $(p, p)$ on complex manifolds. With every analytic subset is associated a current of integration over its set of regular points and all such currents are positive and closed. The important closedness property is proved here via the Skoda-El Mir extension theorem. Positive currents have become an important tool for the study of global geometric problems as well as for questions related to local algebra and intersection theory. We develope here a differential geometric approach to intersection theory through a detailed study of wedge products of closed positive currents (Monge-Ampère operators). The Lelong-Poincaré equation and the Jensen-Lelong formula are basic in this context, providing a useful tool for studying the location and multiplicities of zeroes of entire functions on $\mathbb{C}^{n}$ or on a manifold, in relation with the growth at infinity. Lelong numbers of closed positive currents are then introduced; these numbers can be seen as a generalization to currents of the notion of multiplicity of a germ of analytic set at a singular point. We prove various properties of Lelong numbers (e.g. comparison theorems, semi-continuity theorem of Siu, transformation under holomorphic maps). As an application to Number Theory, we prove a general Schwarz lemma in $\mathbb{C}^{n}$ and derive from it Bombieri's theorem on algebraic values of meromorphic maps and the famous theorems of Gelfond-Schneider and Baker on the transcendence of exponentials and logarithms of algebraic numbers.

## § 1. Basic Concepts of Positivity

## $\S$ 1.A. Positive and Strongly Positive Forms

Let $V$ be a complex vector space of dimension $n$ and $\left(z_{1}, \ldots, z_{n}\right)$ coordinates on $V$. We denote by $\left(\partial / \partial z_{1}, \ldots, \partial / \partial z_{n}\right)$ the corresponding basis of $V$, by $\left(d z_{1}, \ldots, d z_{n}\right)$ its dual basis in $V^{\star}$ and consider the exterior algebra

$$
\Lambda V_{\mathbb{C}}^{\star}=\bigoplus \Lambda^{p, q} V^{\star}, \quad \Lambda^{p, q} V^{\star}=\Lambda^{p} V^{\star} \otimes \Lambda^{q} \overline{V^{\star}}
$$

We are of course especially interested in the case where $V=T_{x} X$ is the tangent space to a complex manifold $X$, but we want to emphasize here that our considerations only involve linear algebra. Let us first observe that $V$ has a canonical orientation, given by the ( $n, n$ )-form

$$
\tau(z)=\mathrm{i} d z_{1} \wedge d \bar{z}_{1} \wedge \ldots \wedge \mathrm{i} d z_{n} \wedge d \bar{z}_{n}=2^{n} d x_{1} \wedge d y_{1} \wedge \ldots \wedge d x_{n} \wedge d y_{n}
$$

where $z_{j}=x_{j}+i y_{j}$. In fact, if $\left(w_{1}, \ldots, w_{n}\right)$ are other coordinates, we find

$$
\begin{aligned}
& d w_{1} \wedge \ldots \wedge d w_{n}=\operatorname{det}\left(\partial w_{j} / \partial z_{k}\right) d z_{1} \wedge \ldots \wedge d z_{n} \\
& \tau(w)=\left|\operatorname{det}\left(\partial w_{j} / \partial z_{k}\right)\right|^{2} \tau(z)
\end{aligned}
$$

In particular, a complex manifold always has a canonical orientation. More generally, natural positivity concepts for ( $p, p$ )-forms can be defined.
(1.1) Definition. $A(p, p)$-form $u \in \Lambda^{p, p} V^{\star}$ is said to be positive if for all $\alpha_{j} \in V^{\star}$, $1 \leqslant j \leqslant q=n-p$, then

$$
u \wedge \mathrm{i} \alpha_{1} \wedge \bar{\alpha}_{1} \wedge \ldots \wedge \mathrm{i} \alpha_{q} \wedge \bar{\alpha}_{q}
$$

is a positive $(n, n)$-form. $A(q, q)$-form $v \in \Lambda^{q, q} V^{\star}$ is said to be strongly positive if $v$ is a convex combination

$$
v=\sum \gamma_{s} \mathrm{i} \alpha_{s, 1} \wedge \bar{\alpha}_{s, 1} \wedge \ldots \wedge \mathrm{i} \alpha_{s, q} \wedge \bar{\alpha}_{s, q}
$$

where $\alpha_{s, j} \in V^{\star}$ and $\gamma_{s} \geqslant 0$.
(1.2) Example. Since $\mathrm{i}^{p}(-1)^{p(p-1) / 2}=\mathrm{i}^{p^{2}}$, we have the commutation rules

$$
\begin{aligned}
& \mathrm{i} \alpha_{1} \wedge \bar{\alpha}_{1} \wedge \ldots \wedge \mathrm{i} \alpha_{p} \wedge \bar{\alpha}_{p}=i^{p^{2}} \alpha \wedge \bar{\alpha}, \quad \forall \alpha=\alpha_{1} \wedge \ldots \wedge \alpha_{p} \in \Lambda^{p, 0} V^{\star}, \\
& i^{p^{2}} \beta \wedge \bar{\beta} \wedge i^{m^{2}} \gamma \wedge \bar{\gamma}=i^{(p+m)^{2}} \beta \wedge \gamma \wedge \overline{\beta \wedge \gamma}, \quad \forall \beta \in \Lambda^{p, 0} V^{\star}, \forall \gamma \in \Lambda^{m, 0} V^{\star} .
\end{aligned}
$$

Take $m=q$ to be the complementary degree of $p$. Then $\beta \wedge \gamma=\lambda d z_{1} \wedge \ldots \wedge d z_{n}$ for some $\lambda \in \mathbb{C}$ and $i^{n^{2}} \beta \wedge \gamma \wedge \overline{\beta \wedge \gamma}=|\lambda|^{2} \tau(z)$. If we set $\gamma=\alpha_{1} \wedge \ldots \wedge \alpha_{q}$, we find that $i^{p^{2}} \beta \wedge \bar{\beta}$ is a positive $(p, p)$-form for every $\beta \in \Lambda^{p, 0} V^{\star}$; in particular, strongly positive forms are positive.

The sets of positive and strongly positive forms are closed convex cones, i.e. closed and stable under convex combinations. By definition, the positive cone is dual to the strongly positive cone via the pairing

$$
\begin{align*}
\Lambda^{p, p} V^{\star} \times \Lambda^{q, q} V^{\star} & \longrightarrow \mathbb{C} \\
(u, v) & \longmapsto u \wedge v / \tau, \tag{1.3}
\end{align*}
$$

that is, $u \in \Lambda^{p, p} V^{\star}$ is positive if and only if $u \wedge v \geqslant 0$ for all strongly positive forms $v \in \Lambda^{q, q} V^{\star}$. Since the bidual of an arbitrary convex cone $\Gamma$ is equal to its closure $\bar{\Gamma}$, we also obtain that $v$ is strongly positive if and only if $v \wedge u=u \wedge v$ is $\geqslant 0$ for all positive forms $u$. Later on, we will need the following elementary lemma.
(1.4) Lemma. Let $\left(z_{1}, \ldots, z_{n}\right)$ be arbitrary coordinates on $V$. Then $\Lambda^{p, p} V^{\star}$ admits a basis consisting of strongly positive forms

$$
\beta_{s}=\mathrm{i} \beta_{s, 1} \wedge \bar{\beta}_{s, 1} \wedge \ldots \wedge \mathrm{i} \beta_{s, p} \wedge \bar{\beta}_{s, p}, \quad 1 \leqslant s \leqslant\binom{ n}{p}^{2}
$$

where each $\beta_{s, l}$ is of the type $d z_{j} \pm d z_{k}$ or $d z_{j} \pm \mathrm{i} d z_{k}, 1 \leqslant j, k \leqslant n$.
Proof. Since one can always extract a basis from a set of generators, it is sufficient to see that the family of forms of the above type generates $\Lambda^{p, p} V^{\star}$. This follows from the identities

$$
\begin{aligned}
4 d z_{j} \wedge d \bar{z}_{k}= & \left(d z_{j}+d z_{k}\right) \wedge\left(\overline{d z_{j}+d z_{k}}\right)-\left(d z_{j}-d z_{k}\right) \wedge\left(\overline{d z_{j}-d z_{k}}\right) \\
& +\mathrm{i}\left(d z_{j}+\mathrm{i} d z_{k}\right) \wedge\left(\overline{d z_{j}+\mathrm{i} d z_{k}}\right)-\mathrm{i}\left(d z_{j}-\mathrm{i} d z_{k}\right) \wedge\left(\overline{d z_{j}-\mathrm{i} d z_{k}}\right),
\end{aligned}
$$

$$
d z_{j_{1}} \wedge \ldots \wedge d z_{j_{p}} \wedge d \bar{z}_{k_{1}} \wedge \ldots \wedge d \bar{z}_{k_{p}}= \pm \bigwedge_{1 \leqslant s \leqslant p} d z_{j_{s}} \wedge d \bar{z}_{k_{s}}
$$

(1.5) Corollary. All positive forms $u$ are real, i.e. satisfy $\bar{u}=u$. In terms of coordinates, if $u=i^{p^{2}} \sum_{|I|=|J|=p} u_{I, J} d z_{I} \wedge d \bar{z}_{J}$, then the coefficients satisfy the hermitian symmetry relation $\overline{u_{I, J}}=u_{J, I}$.

Proof. Clearly, every strongly positive $(q, q)$-form is real. By Lemma 1.4, these forms generate over $\mathbb{R}$ the real elements of $\Lambda^{q, q} V^{\star}$, so we conclude by duality that positive $(p, p)$-forms are also real.
(1.6) Criterion. A form $u \in \Lambda^{p, p} V^{\star}$ is positive if and only if its restriction $u_{\mid S}$ to every p-dimensional subspace $S \subset V$ is a positive volume form on $S$.

Proof. If $S$ is an arbitrary $p$-dimensional subspace of $V$ we can find linear coordinates $\left(z_{1}, \ldots, z_{n}\right)$ on $V$ such that $S=\left\{z_{p+1}=\ldots=z_{n}=0\right\}$. Then

$$
u_{\upharpoonright S}=\lambda_{S} \mathrm{i} d z_{1} \wedge d \bar{z}_{1} \wedge \ldots \wedge \mathrm{i} d z_{p} \wedge d \bar{z}_{p}
$$

where $\lambda_{S}$ is given by

$$
u \wedge \mathrm{i} d z_{p+1} \wedge d \bar{z}_{p+1} \wedge \ldots \wedge \mathrm{i} d z_{n} \wedge d \bar{z}_{n}=\lambda_{S} \tau(z)
$$

If $u$ is positive then $\lambda_{S} \geqslant 0$ so $u_{\upharpoonright S}$ is positive for every $S$. The converse is true because the ( $n-p, n-p$ )-forms $\bigwedge_{j>p} \mathrm{i} d z_{j} \wedge d \bar{z}_{j}$ generate all strongly positive forms when $S$ runs over all $p$-dimensional subspaces.
(1.7) Corollary. $A$ form $u=\mathrm{i} \sum_{j, k} u_{j k} d z_{j} \wedge d \bar{z}_{k}$ of bidegree $(1,1)$ is positive if and only if $\xi \mapsto \sum u_{j k} \xi_{j} \bar{\xi}_{k}$ is a semi-positive hermitian form on $\mathbb{C}^{n}$.

Proof. If $S$ is the complex line generated by $\xi$ and $t \mapsto t \xi$ the parametrization of $S$, then $u_{\upharpoonright S}=\left(\sum u_{j k} \xi_{j} \bar{\xi}_{k}\right) \mathrm{i} d t \wedge d \bar{t}$.

Observe that there is a canonical one-to-one correspondence between hermitian forms and real ( 1,1 )-forms on $V$. The correspondence is given by

$$
\begin{equation*}
h=\sum_{1 \leqslant j, k \leqslant n} h_{j k}(z) d z_{j} \otimes d \bar{z}_{k} \longmapsto u=\mathrm{i} \sum_{1 \leqslant j, k \leqslant n} h_{j k}(z) d z_{j} \wedge d \bar{z}_{k} \tag{1.8}
\end{equation*}
$$

and does not depend on the choice of coordinates: indeed, as $\bar{h}_{j k}=h_{k j}$, one finds immediately

$$
u(\xi, \eta)=\mathrm{i} \sum h_{j k}(z)\left(\xi_{j} \bar{\eta}_{k}-\eta_{j} \bar{\xi}_{k}\right)=-2 \operatorname{Im} h(\xi, \eta), \quad \forall \xi, \eta \in T X
$$

Moreover, $h$ is $\geqslant 0$ as a hermitian form if and only if $u \geqslant 0$ as a ( 1,1 )-form. A diagonalization of $h$ shows that every positive ( 1,1 )-form $u \in \Lambda^{1,1} V^{\star}$ can be written

$$
u=\sum_{1 \leqslant j \leqslant r} \mathrm{i} \gamma_{j} \wedge \bar{\gamma}_{j}, \quad \gamma \in V^{\star}, \quad r=\text { rank of } u
$$

in particular, every positive ( 1,1 )-form is strongly positive. By duality, this is also true for ( $n-1, n-1$ )-forms.
(1.9) Corollary. The notions of positive and strongly positive $(p, p)$-forms coincide for $p=0,1, n-1, n$.
(1.10) Remark. It is not difficult to see, however, that positivity and strong positivity differ in all bidegrees $(p, p)$ such that $2 \leqslant p \leqslant n-2$. Indeed, a positive form $i^{p^{2}} \beta \wedge \bar{\beta}$ with $\beta \in \Lambda^{p, 0} V^{\star}$ is strongly positive if and only if $\beta$ is decomposable as a product $\beta_{1} \wedge \ldots \wedge \beta_{p}$. To see this, suppose that

$$
i^{p^{2}} \beta \wedge \bar{\beta}=\sum_{1 \leqslant j \leqslant N} i^{p^{2}} \gamma_{j} \wedge \bar{\gamma}_{j}
$$

where all $\gamma_{j} \in \Lambda^{p, 0} V^{\star}$ are decomposable. Take $N$ minimal. The equality can be also written as an equality of hermitian forms $|\beta|^{2}=\sum\left|\gamma_{j}\right|^{2}$ if $\beta, \gamma_{j}$ are seen as linear forms on $\Lambda^{p} V$. The hermitian form $|\beta|^{2}$ has rank one, so we must have $N=1$ and $\beta=\lambda \gamma_{j}$, as desired. Note that there are many non decomposable $p$-forms in all degrees $p$ such that $2 \leqslant p \leqslant n-2$, e.g. $\left(d z_{1} \wedge d z_{2}+d z_{3} \wedge d z_{4}\right) \wedge d z_{5} \wedge \ldots \wedge d z_{p+2}:$ if a $p$-form is decomposable, the vector space of its contractions by elements of $\bigwedge^{p-1} V$ is a $p$-dimensional subspace of $V^{\star}$; in the above example the dimension is $p+2$.
(1.11) Proposition. If $u_{1}, \ldots, u_{s}$ are positive forms, all of them strongly positive (resp. all except perhaps one), then $u_{1} \wedge \ldots \wedge u_{\text {s }}$ is strongly positive (resp. positive).

Proof. Immediate consequence of Def. 1.1. Observe however that the wedge product of two positive forms is not positive in general (otherwise we would infer that positivity coincides with strong positivity).
(1.12) Proposition. If $\Phi: W \longrightarrow V$ is a complex linear map and $u \in \Lambda^{p, p} V^{\star}$ is (strongly) positive, then $\Phi^{\star} u \in \Lambda^{p, p} W^{\star}$ is (strongly) positive.

Proof. This is clear for strong positivity, since

$$
\Phi^{\star}\left(\mathrm{i} \alpha_{1} \wedge \bar{\alpha}_{1} \wedge \ldots \wedge \mathrm{i} \alpha_{p} \wedge \bar{\alpha}_{p}\right)=\mathrm{i} \beta_{1} \wedge \bar{\beta}_{1} \wedge \ldots \wedge \mathrm{i} \beta_{p} \wedge \bar{\beta}_{p}
$$

with $\beta_{j}=\Phi^{\star} \alpha_{j} \in W^{\star}$, for all $\alpha_{j} \in V^{\star}$. For $u$ positive, we may apply Criterion 1.6: if $S$ is a $p$-dimensional subspace of $W$, then $u_{\mid \Phi(S)}$ and $\left(\Phi^{\star} u\right)_{\mid S}=\left(\Phi_{\uparrow} S\right)^{\star} u_{\mid \Phi(S)}$ are positive when $\Phi_{\mid S}: S \longrightarrow \Phi(S)$ is an isomorphism; otherwise we get $\left(\Phi^{\star} u\right)_{\upharpoonright S}=0$.

## $\S$ 1.B. Positive Currents

The duality between the positive and strongly positive cones of forms can be used to define corresponding positivity notions for currents.
(1.13) Definition. A current $T \in \mathscr{D}_{p, p}^{\prime}(X)$ is said to be positive (resp. strongly positive) if $\langle T, u\rangle \geqslant 0$ for all test forms $u \in \mathscr{D}_{p, p}(X)$ that are strongly positive (resp. positive) at each point. The set of positive (resp. strongly positive) currents of bidimension ( $p, p$ ) will be denoted

$$
\mathscr{D}_{p, p}^{\prime+}(X), \quad \text { resp. } \mathscr{D}_{p, p}^{\prime \oplus}(X)
$$

It is clear that (strong) positivity is a local property and that the sets $\mathscr{D}_{p, p}^{\prime \oplus}(X) \subset$ $\mathscr{D}_{p, p}^{\prime+}(X)$ are closed convex cones with respect to the weak topology. Another way of stating Def. 1.13 is:
$T$ is positive (strongly positive) if and only if $T \wedge u \in \mathscr{D}_{0,0}^{\prime}(X)$ is a positive measure for all strongly positive (positive) forms $u \in \mathscr{C}_{p, p}^{\infty}(X)$.

This is so because a distribution $S \in \mathscr{D}^{\prime}(X)$ such that $S(f) \geqslant 0$ for every non-negative function $f \in \mathscr{D}(X)$ is a positive measure.
(1.14) Proposition. Every positive current $T=i^{(n-p)^{2}} \sum T_{I, J} d z_{I} \wedge d \bar{z}_{J}$ in $\mathscr{D}_{p, p}^{\prime+}(X)$ is real and of order 0 , i.e. its coefficients $T_{I, J}$ are complex measures and satisfy $\overline{T_{I, J}}=T_{J, I}$ for all multi-indices $|I|=|J|=n-p$. Moreover $T_{I, I} \geqslant 0$, and the absolute values $\left|T_{I, J}\right|$ of the measures $T_{I, J}$ satisfy the inequality

$$
\lambda_{I} \lambda_{J}\left|T_{I, J}\right| \leqslant 2^{p} \sum_{M} \lambda_{M}^{2} T_{M, M}, \quad I \cap J \subset M \subset I \cup J
$$

where $\lambda_{k} \geqslant 0$ are arbitrary coefficients and $\lambda_{I}=\prod_{k \in I} \lambda_{k}$.
Proof. Since positive forms are real, positive currents have to be real by duality. Let us denote by $K=\complement I$ and $L=\complement J$ the ordered complementary multi-indices of $I, J$ in $\{1,2, \ldots, n\}$. The distribution $T_{I, I}$ is a positive measure because

$$
T_{I, I} \tau=T \wedge i^{p^{2}} d z_{K} \wedge d \bar{z}_{K} \geqslant 0
$$

On the other hand, the proof of Lemma 1.4 yields

$$
\begin{aligned}
T_{I, J} \tau & = \pm T \wedge i^{p^{2}} d z_{K} \wedge d \bar{z}_{L}=\sum_{a \in(\mathbb{Z} / 4 \mathbb{Z})^{p}} \varepsilon_{a} T \wedge \gamma_{a} \quad \text { where } \\
\gamma_{a} & =\bigwedge_{1 \leqslant s \leqslant p} \frac{\mathrm{i}}{4}\left(d z_{k_{s}}+i^{a_{s}} d z_{l_{s}}\right) \wedge\left(\overline{d z_{k_{s}}+i^{a_{s}} d z_{l_{s}}}\right), \quad \varepsilon_{a}= \pm 1, \pm i
\end{aligned}
$$

Now, each $T \wedge \gamma_{a}$ is a positive measure, hence $T_{I, J}$ is a complex measure and

$$
\begin{aligned}
\left|T_{I, J}\right| \tau & \leqslant \sum_{a} T \wedge \gamma_{a}=T \wedge \sum_{a} \gamma_{a} \\
& =T \wedge \bigwedge_{1 \leqslant s \leqslant p}\left(\sum_{a_{s} \in \mathbb{Z} / 4 \mathbb{Z}} \frac{\mathrm{i}}{4}\left(d z_{k_{s}}+i^{a_{s}} d z_{l_{s}}\right) \wedge\left(\overline{d z_{k_{s}}+i^{a_{s}} d z_{l_{s}}}\right)\right) \\
& =T \wedge \bigwedge_{1 \leqslant s \leqslant p}\left(\mathrm{i} d z_{k_{s}} \wedge d \bar{z}_{k_{s}}+\mathrm{i} d z_{l_{s}} \wedge d \bar{z}_{l_{s}}\right)
\end{aligned}
$$

The last wedge product is a sum of at most $2^{p}$ terms, each of which is of the type $i^{p^{2}} d z_{M} \wedge d \bar{z}_{M}$ with $|M|=p$ and $M \subset K \cup L$. Since $T \wedge i^{p^{2}} d z_{M} \wedge d \bar{z}_{M}=T_{\mathrm{C} M, \mathrm{C} M} \tau$ and $\complement M \supset \complement K \cap \complement L=I \cap J$, we find

$$
\left|T_{I, J}\right| \leqslant 2^{p} \sum_{M \supset I \cap J} T_{M, M}
$$

Now, consider a change of coordinates $\left(z_{1}, \ldots, z_{n}\right)=\Lambda w=\left(\lambda_{1} w_{1}, \ldots, \lambda_{n} w_{n}\right)$ with $\lambda_{1}, \ldots, \lambda_{n}>0$. In the new coordinates, the current $T$ becomes $\Lambda^{\star} T$ and its coefficients become $\lambda_{I} \lambda_{J} T_{I, J}(\Lambda w)$. Hence, the above inequality implies

$$
\lambda_{I} \lambda_{J}\left|T_{I, J}\right| \leqslant 2^{p} \sum_{M \supset I \cap J} \lambda_{M}^{2} T_{M, M}
$$

This inequality is still true for $\lambda_{k} \geqslant 0$ by passing to the limit. The inequality of Prop. 1.14 follows when all coefficients $\lambda_{k}, k \notin I \cup J$, are replaced by 0 , so that $\lambda_{M}=0$ for $M \not \subset I \cup J$.
(1.15) Remark. If $T$ is of order 0 , we define the mass measure of $T$ by $\|T\|=\sum\left|T_{I, J}\right|$ (of course $\|T\|$ depends on the choice of coordinates). By the Radon-Nikodym theorem, we can write $T_{I, J}=f_{I, J}\|T\|$ with a Borel function $f_{I, J}$ such that $\sum\left|f_{I, J}\right|=1$. Hence

$$
T=\|T\| f, \quad \text { where } \quad f=i^{(n-p)^{2}} \sum f_{I, J} d z_{I} \wedge d \bar{z}_{J}
$$

Then $T$ is (strongly) positive if and only if the form $f(x) \in \Lambda^{n-p, n-p} T_{x}^{\star} X$ is (strongly) positive at $\|T\|$-almost all points $x \in X$. Indeed, this condition is clearly sufficient. On the other hand, if $T$ is (strongly) positive and $u_{j}$ is a sequence of forms with constant coefficients in $\Lambda^{p, p} T^{\star} X$ which is dense in the set of strongly positive (positive) forms, then $T \wedge u_{j}=\|T\| f \wedge u_{j}$, so $f(x) \wedge u_{j}$ has to be a positive ( $n, n$ )-form except perhaps for $x$ in a set $N\left(u_{j}\right)$ of $\|T\|$-measure 0 . By a simple density argument, we see that $f(x)$ is (strongly) positive outside the $\|T\|$-negligible set $N=\bigcup N\left(u_{j}\right)$.

As a consequence of this proof, $T$ is positive (strongly positive) if and only if $T \wedge u$ is a positive measure for all strongly positive (positive) forms $u$ of bidegree ( $p, p$ ) with constant coefficients in the given coordinates $\left(z_{1}, \ldots, z_{n}\right)$. It follows that if $T$ is (strongly) positive in a coordinate patch $\Omega$, then the convolution $T \star \rho_{\varepsilon}$ is (strongly) positive in $\Omega_{\varepsilon}=\{x \in \Omega ; d(x, \partial \Omega)>\varepsilon\}$.
(1.16) Corollary. If $T \in \mathscr{D}_{p, p}^{\prime}(X)$ and $v \in C_{s, s}^{0}(X)$ are positive, one of them (resp. both of them) strongly positive, then the wedge product $T \wedge v$ is a positive (resp. strongly positive) current.

This follows immediately from Remark 1.15 and Prop. 1.11 for forms. Similarly, Prop. 1.12 on pull-backs of positive forms easily shows that positivity properties of currents are preserved under direct or inverse images by holomorphic maps.
(1.17) Proposition. Let $\Phi: X \longrightarrow Y$ be a holomorphic map between complex analytic manifolds.
a) If $T \in \mathscr{D}_{p, p}^{\prime+}(X)$ and $\Phi_{\mid \text {Supp } T}$ is proper, then $\Phi_{\star} T \in \mathscr{D}_{p, p}^{\prime+}(Y)$.
b) If $T \in \mathscr{D}_{p, p}^{\prime+}(Y)$ and if $\Phi$ is a submersion with $m$-dimensional fibers, then $\Phi^{\star} T \in$ $\mathscr{D}_{p+m, p+m}^{\prime+}(X)$.

Similar properties hold for strongly positive currents.

## § 1.C. Basic Examples of Positive Currents

We present here two fundamental examples which will be of interest in many circumstances.
$\S$ 1.18. Current Associated to a Plurisubharmonic Function. Let $X$ be a complex manifold and $u \in \operatorname{Psh}(X) \cap L_{\text {loc }}^{1}(X)$ a plurisubharmonic function. Then

$$
T=\mathrm{i} d^{\prime} d^{\prime \prime} u=\mathrm{i} \sum_{1 \leqslant j, k \leqslant n} \frac{\partial^{2} u}{\partial z_{j} \partial \bar{z}_{k}} d z_{j} \wedge d \bar{z}_{k}
$$

is a positive current of bidegree $(1,1)$. Moreover $T$ is closed (we always mean here $d$ closed, that is, $d T=0$ ). Assume conversely that $\Theta$ is a closed real $(1,1)$-current on $X$. Poincaré's lemma implies that every point $x_{0} \in X$ has a neighborhood $\Omega_{0}$ such that $\Theta=d S$ with $S \in \mathscr{D}_{1}^{\prime}\left(\Omega_{0}, \mathbb{R}\right)$. Write $S=S^{1,0}+S^{0,1}$, where $S^{0,1}=\overline{S^{1,0}}$. Then $d^{\prime \prime} S=\Theta^{0,2}=0$, and the Dolbeault-Grothendieck lemma shows that $S^{0,1}=d^{\prime \prime} v$ on some neighborhood $\Omega \subset \Omega_{0}$, with $v \in \mathscr{D}^{\prime}(\Omega, \mathbb{C})$. Thus

$$
\begin{aligned}
& S=\overline{d^{\prime \prime} v}+d^{\prime \prime} v=d^{\prime} \bar{v}+d^{\prime \prime} v \\
& \Theta=d S=d^{\prime} d^{\prime \prime}(v-\bar{v})=\mathrm{i} d^{\prime} d^{\prime \prime} u
\end{aligned}
$$

where $u=2 \operatorname{Re} v \in \mathscr{D}^{\prime}(\Omega, \mathbb{R})$. If $\Theta \in \mathscr{C}_{1,1}^{\infty}(X)$, the hypoellipticity of $d^{\prime \prime}$ in bidegree $(p, 0)$ shows that $d^{\prime} u$ is of class $\mathscr{C}^{\infty}$, so $u \in \mathscr{C}^{\infty}(\Omega)$. When $\Theta$ is positive, the distribution $u$ is a plurisubharmonic function (Th. I.3.31). We have thus proved:
(1.19) Proposition. If $\Theta \in \mathscr{D}_{n-1, n-1}^{+}(X)$ is a closed positive current of bidegree $(1,1)$, then for every point $x_{0} \in X$ there exists a neighborhood $\Omega$ of $x_{0}$ and $u \in \operatorname{Psh}(\Omega)$ such that $\Theta=\mathrm{i} d^{\prime} d^{\prime \prime} u$.
§ 1.20. Current of Integration on a Complex Submanifold. Let $Z \subset X$ be a closed $p$-dimensional complex submanifold with its canonical orientation and $T=[Z]$. Then $T \in \mathscr{D}_{p, p}^{\prime}(X)$. Indeed, every $(r, s)$-form of total degree $r+s=2 p$ has zero restriction to $Z$ unless $(r, s)=(p, p)$, therefore we have $[Z] \in \mathscr{D}_{p, p}^{\prime}(X)$. Now, if $u \in \mathscr{D}_{p, p}(X)$ is a positive test form, then $u_{\mid Z}$ is a positive volume form on $Z$ by Criterion 1.6, therefore

$$
\langle[Z], u\rangle=\int_{Z} u_{\curlyvee Z} \geqslant 0
$$

In this example the current $[Z]$ is also closed, because $d[Z]= \pm[\partial Z]=0$ by Stokes' theorem.

## §1.D. Trace Measure and Wirtinger's Inequality

We discuss now some questions related to the concept of area on complex submanifolds. Assume that $X$ is equipped with a hermitian metric $h$, i.e. a positive definite hermitian form $h=\sum h_{j k} d z_{j} \otimes d \bar{z}_{k}$ of class $\mathscr{C}^{\infty}$; we denote by $\omega=\mathrm{i} \sum h_{j k} d z_{j} \wedge d \bar{z}_{k} \in \mathscr{C}_{1,1}^{\infty}(X)$ the associated positive (1,1)-form.
(1.21) Definition. For every $T \in \mathscr{D}_{p, p}^{\prime+}(X)$, the trace measure of $T$ with respect to $\omega$ is the positive measure

$$
\sigma_{T}=\frac{1}{2^{p} p!} T \wedge \omega^{p} .
$$

If $\left(\zeta_{1}, \ldots, \zeta_{n}\right)$ is an orthonormal frame of $T^{\star} X$ with respect to $h$ on an open subset $U \subset X$, we may write

$$
\begin{aligned}
& \omega=\mathrm{i} \sum_{1 \leqslant j \leqslant n} \zeta_{j} \wedge \bar{\zeta}_{j}, \quad \omega^{p}=i^{p^{2}} p!\sum_{|K|=p} \zeta_{K} \wedge \bar{\zeta}_{K}, \\
& T=i^{(n-p)^{2}} \sum_{|I|=|J|=n-p} T_{I, J} \zeta_{I} \wedge \bar{\zeta}_{J}, \quad T_{I, J} \in \mathscr{D}^{\prime}(U),
\end{aligned}
$$

where $\zeta_{I}=\zeta_{i_{1}} \wedge \ldots \wedge \zeta_{i_{n-p}}$. An easy computation yields

$$
\begin{equation*}
\sigma_{T}=2^{-p}\left(\sum_{|I|=n-p} T_{I, I}\right) \mathrm{i} \zeta_{1} \wedge \bar{\zeta}_{1} \wedge \ldots \wedge \mathrm{i} \zeta_{n} \wedge \bar{\zeta}_{n} \tag{1.22}
\end{equation*}
$$

For $X=\mathbb{C}^{n}$ with the standard hermitian metric $h=\sum d z_{j} \otimes d \bar{z}_{j}$, we get in particular

$$
\sigma_{T}=2^{-p}\left(\sum_{|I|=n-p} T_{I, I}\right) \mathrm{i} d z_{1} \wedge d \bar{z}_{1} \wedge \ldots \wedge \mathrm{i} d z_{n} \wedge d \bar{z}_{n} .
$$

Proposition 1.14 shows that the mass measure $\|T\|=\sum\left|T_{I, J}\right|$ of a positive current $T$ is always dominated by $C \sigma_{T}$ where $C>0$ is a constant. It follows easily that the weak topology of $\mathscr{D}_{p}^{\prime}(X)$ and of $\mathscr{D}_{p}^{0 \prime}(X)$ coincide on $\mathscr{D}_{p}^{\prime+}(X)$, which is moreover a metrizable subspace: its weak topology is in fact defined by the collection of semi-norms $T \longmapsto$ $\left|\left\langle T, f_{\nu}\right\rangle\right|$ where $\left(f_{\nu}\right)$ is an arbitrary dense sequence in $\mathscr{D}_{p}(X)$. By the Banach-Alaoglu theorem, the unit ball in the dual of a Banach space is weakly compact, thus:
(1.23) Proposition. Let $\delta$ be a positive continuous function on $X$. Then the set of currents $T \in \mathscr{D}_{p}^{\prime+}(X)$ such that $\int_{X} \delta T \wedge \omega^{p} \leqslant 1$ is weakly compact.

Proof. Note that our set is weakly closed, since a weak limit of positive currents is positive and $\int_{X} \delta T \wedge \omega^{p}=\sup \left\langle T, \theta \delta \omega^{p}\right\rangle$ when $\theta$ runs over all elements of $\mathscr{D}(X)$ such that $0 \leqslant \theta \leqslant 1$.

Now, let $Z$ be a $p$-dimensional complex analytic submanifold of $X$. We claim that

$$
\begin{equation*}
\sigma_{[Z]}=\frac{1}{2^{p} p!}[Z] \wedge \omega^{p}=\text { Riemannian volume measure on } Z . \tag{1.24}
\end{equation*}
$$

This result is in fact a special case of the following important inequality.
(1.25) Wirtinger's inequality. Let $Y$ be an oriented real submanifold of class $C^{1}$ and dimension $2 p$ in $X$, and let $d V_{Y}$ be the Riemannian volume form on $Y$ associated with the metric $h_{\lceil Y}$. Set

$$
\frac{1}{2^{p} p!} \omega_{\upharpoonright Y}^{p}=\alpha d V_{Y}, \quad \alpha \in C^{0}(Y)
$$

Then $|\alpha| \leqslant 1$ and the equality holds if and only if $Y$ is a complex analytic submanifold of $X$. In that case $\alpha=1$ if the orientation of $Y$ is the canonical one, $\alpha=-1$ otherwise.

Proof. The restriction $\omega_{\lceil Y}$ is a real antisymmetric 2-form on $T Y$. At any point $z \in Y$, we can thus find an oriented orthonormal $\mathbb{R}$-basis $\left(e_{1}, e_{2}, \ldots, e_{2 p}\right)$ of $T_{z} Y$ such that

$$
\begin{aligned}
\frac{1}{2} \omega & =\sum_{1 \leqslant k \leqslant p} \alpha_{k} e_{2 k-1}^{\star} \wedge e_{2 k}^{\star} \quad \text { on } T_{z} Y, \quad \text { where } \\
\alpha_{k} & =\frac{1}{2} \omega\left(e_{2 k-1}, e_{2 k}\right)=-\operatorname{Im} h\left(e_{2 k-1}, e_{2 k}\right)
\end{aligned}
$$

We have $d V_{Y}=e_{1}^{\star} \wedge \ldots \wedge e_{2 p}^{\star}$ by definition of the Riemannian volume form. By taking the $p$-th power of $\omega$, we get

$$
\frac{1}{2^{p} p!} \omega_{\upharpoonright T_{z} Y}^{p}=\alpha_{1} \ldots \alpha_{p} e_{1}^{\star} \wedge \ldots \wedge e_{2 p}^{\star}=\alpha_{1} \ldots \alpha_{p} d V_{Y}
$$

Since $\left(e_{k}\right)$ is an orthonormal $\mathbb{R}$-basis, we have $\operatorname{Re} h\left(e_{2 k-1}, e_{2 k}\right)=0$, thus $h\left(e_{2 k-1}, e_{2 k}\right)=$ $-\mathrm{i} \alpha_{k}$. As $\left|e_{2 k-1}\right|=\left|e_{2 k}\right|=1$, we get

$$
0 \leqslant\left|e_{2 k} \pm J e_{2 k-1}\right|^{2}=2\left(1 \pm \operatorname{Re} h\left(J e_{2 k-1}, e_{2 k}\right)\right)=2\left(1 \pm \alpha_{k}\right)
$$

Therefore

$$
\left|\alpha_{k}\right| \leqslant 1, \quad|\alpha|=\left|\alpha_{1} \ldots \alpha_{p}\right| \leqslant 1,
$$

with equality if and only if $\alpha_{k}= \pm 1$ for all $k$, i.e. $e_{2 k}= \pm J e_{2 k-1}$. In this case $T_{z} Y \subset T_{z} X$ is a complex vector subspace at every point $z \in Y$, thus $Y$ is complex analytic by Lemma I.4.23. Conversely, assume that $Y$ is a $\mathbb{C}$-analytic submanifold and that $\left(e_{1}, e_{3}, \ldots, e_{2 p-1}\right)$ is an orthonormal complex basis of $T_{z} Y$. If $e_{2 k}:=J e_{2 k-1}$, then $\left(e_{1}, \ldots, e_{2 p}\right)$ is an orthonormal $\mathbb{R}$-basis corresponding to the canonical orientation and

$$
\frac{1}{2} \omega_{\lceil Y}=\sum_{1 \leqslant k \leqslant p} e_{2 k-1}^{\star} \wedge e_{2 k}^{\star}, \quad \frac{1}{2^{p} p!} \omega_{\Gamma Y}^{p}=e_{1}^{\star} \wedge \ldots \wedge e_{2 p}^{\star}=d V_{Y}
$$

Note that in the case of the standard hermitian metric $\omega$ on $X=\mathbb{C}^{n}$, the form $\omega=$ $\mathrm{i} \sum d z_{j} \wedge d \bar{z}_{j}=d\left(\mathrm{i} \sum z_{j} d \bar{z}_{j}\right)$ is globally exact. Under this hypothesis, we are going to see that $\mathbb{C}$-analytic submanifolds are always minimal surfaces for the Plateau problem, which consists in finding a compact subvariety $Y$ of minimal area with prescribed boundary $\partial Y$.
(1.26) Theorem. Assume that the $(1,1)$-form $\omega$ is exact, say $\omega=d \gamma$ with $\gamma \in$ $\mathscr{C}_{1}^{\infty}(X, \mathbb{R})$, and let $Y, Z \subset X$ be $(2 p)$-dimensional oriented compact real submanifolds of class $C^{1}$ with boundary. If $\partial Y=\partial Z$ and $Z$ is complex analytic, then

$$
\operatorname{Vol}(Y) \geqslant \operatorname{Vol}(\mathrm{Z})
$$

Proof. Write $\omega=d \gamma$. Wirtinger's inequality and Stokes' theorem imply

$$
\begin{aligned}
& \operatorname{Vol}(Y) \geqslant \frac{1}{2^{p} p!}\left|\int_{Y} \omega^{p}\right|=\frac{1}{2^{p} p!}\left|\int_{Y} d\left(\omega^{p-1} \wedge \gamma\right)\right|=\frac{1}{2^{p} p!}\left|\int_{\partial Y} \omega^{p-1} \wedge \gamma\right| \\
& \operatorname{Vol}(Z)=\frac{1}{2^{p} p!} \int_{Z} \omega^{p}=\frac{1}{2^{p} p!} \int_{\partial Z} \omega^{p-1} \wedge \gamma= \pm \frac{1}{2^{p} p!} \int_{\partial Y} \omega^{p-1} \wedge \gamma
\end{aligned}
$$

## § 2. Closed Positive Currents

## $\S$ 2.A. The Skoda-El Mir Extension Theorem

We first prove the Skoda-El Mir extension theorem ([Skoda 1982], [El Mir 1984]), which shows in particular that a closed positive current defined in the complement of an analytic set $E$ can be extended through $E$ if (and only if) the mass of the current is locally finite near $E$. El Mir simplified Skoda's argument and showed that it is enough to assume $E$ complete pluripolar. We follow here the exposition of Sibony's survey article [Sibony 1985].
(2.1) Definition. $A$ subset $E \subset X$ is said to be complete pluripolar in $X$ if for every point $x_{0} \in X$ there exist a neighborhood $\Omega$ of $x_{0}$ and a function $u \in \operatorname{Psh}(\Omega) \cap L_{\mathrm{loc}}^{1}(\Omega)$ such that $E \cap \Omega=\{z \in \Omega ; u(z)=-\infty\}$.

Note that any closed analytic subset $A \subset X$ is complete pluripolar: if $g_{1}=\ldots=$ $g_{N}=0$ are holomorphic equations of $A$ on an open set $\Omega \subset X$, we can take $u=$ $\log \left(\left|g_{1}\right|^{2}+\ldots+\left|g_{N}\right|^{2}\right)$.
(2.2) Lemma. Let $E \subset X$ be a closed complete pluripolar set. If $x_{0} \in X$ and $\Omega$ is a sufficiently small neighborhood of $x_{0}$, there exists:
a) a function $v \in \operatorname{Psh}(\Omega) \cap \mathscr{C}^{\infty}(\Omega \backslash E)$ such that $v=-\infty$ on $E \cap \Omega$;
b) an increasing sequence $v_{k} \in \operatorname{Psh}(\Omega) \cap \mathscr{C}^{\infty}(\Omega), 0 \leqslant v_{k} \leqslant 1$, converging uniformly to 1 on every compact subset of $\Omega \backslash E$, such that $v_{k}=0$ on a neighborhood of $E \cap \Omega$.

Proof. Assume that $\Omega_{0} \subset \subset X$ is a coordinate patch of $X$ containing $x_{0}$ and that $E \cap \Omega_{0}=$ $\left\{z \in \Omega_{0} ; u(z)=-\infty\right\}, u \in \operatorname{Psh}\left(\Omega_{0}\right)$. In addition, we can achieve $u \leqslant 0$ by shrinking $\Omega_{0}$ and subtracting a constant to $u$. Select a convex increasing function $\chi \in \mathscr{C}^{\infty}([0,1], \mathbb{R})$ such that $\chi(t)=0$ on $[0,1 / 2]$ and $\chi(1)=1$. We set

$$
u_{k}=\chi(\exp (u / k))
$$

Then $0 \leqslant u_{k} \leqslant 1, u_{k}$ is plurisubharmonic on $\Omega_{0}, u_{k}=0$ in a neighborhood $\omega_{k}$ of $E \cap \Omega_{0}$ and $\lim u_{k}=1$ on $\Omega_{0} \backslash E$. Let $\Omega \subset \subset \Omega_{0}$ be a neighborhood of $x_{0}$, let $\delta_{0}=d\left(\Omega, \subset \Omega_{0}\right)$ and $\left.\varepsilon_{k} \in\right] 0, \delta_{0}\left[\right.$ be such that $\varepsilon_{k}<d\left(E \cap \bar{\Omega}, \bar{\Omega} \backslash \omega_{k}\right)$. Then

$$
w_{k}:=\max _{j \leqslant k}\left\{u_{j} \star \rho_{\varepsilon_{j}}\right\} \in \operatorname{Psh}(\Omega) \cap C^{0}(\Omega),
$$

$0 \leqslant w_{k} \leqslant 1, w_{k}=0$ on a neighborhood of $E \cap \Omega$ and $w_{k}$ is an increasing sequence converging to 1 on $\Omega \backslash E$ (note that $w_{k} \geqslant u_{k}$ ). Hence, the convergence is uniform on every compact subset of $\Omega \backslash E$ by Dini's lemma. We may therefore choose a subsequence $w_{k_{s}}$ such that $w_{k_{s}}(z) \geqslant 1-2^{-s}$ on an increasing sequence of open sets $\Omega_{s}^{\prime}$ with $\bigcup \Omega_{s}^{\prime}=\Omega \backslash E$. Then

$$
w(z):=|z|^{2}+\sum_{s=0}^{+\infty}\left(w_{k_{s}}(z)-1\right)
$$

is a strictly plurisubharmonic function on $\Omega$ that is continuous on $\Omega \backslash E$, and $w=-\infty$ on $E \cap \Omega$. Richberg's theorem I.3.40 applied on $\Omega \backslash E$ produces $v \in \operatorname{Psh}(\Omega \backslash E) \cap \mathscr{C}^{\infty}(\Omega \backslash E)$
such that $w \leqslant v \leqslant w+1$. If we set $v=-\infty$ on $E \cap \Omega$, then $v$ is plurisubharmonic on $\Omega$ and has the properties required in a). After subtraction of a constant, we may assume $v \leqslant 0$ on $\Omega$. Then the sequence $\left(v_{k}\right)$ of statement b$)$ is obtained by letting $v_{k}=\chi(\exp (v / k))$.
(2.3) Theorem (El Mir). Let $E \subset X$ be a closed complete pluripolar set and $T \in$ $\mathscr{D}_{p, p}^{\prime+}(X \backslash E)$ a closed positive current. Assume that $T$ has finite mass in a neighborhood of every point of $E$. Then the trivial extension $\tilde{T} \in \mathscr{D}_{p, p}^{\prime+}(X)$ obtained by extending the measures $T_{I, J}$ by 0 on $E$ is closed on $X$.

Proof. The statement is local on $X$, so we may work on a small open set $\Omega$ such that there exists a sequence $v_{k} \in \operatorname{Psh}(\Omega) \cap \mathscr{C}^{\infty}(\Omega)$ as in 2.2 b). Let $\theta \in \mathscr{C}^{\infty}([0,1])$ be a function such that $\theta=0$ on $[0,1 / 3], \theta=1$ on $[2 / 3,1]$ and $0 \leqslant \theta \leqslant 1$. Then $\theta \circ v_{k}=0$ near $E \cap \Omega$ and $\theta \circ v_{k}=1$ for $k$ large on every fixed compact subset of $\Omega \backslash E$. Therefore $\tilde{T}=\lim _{k \rightarrow+\infty}\left(\theta \circ v_{k}\right) T$ and

$$
d^{\prime} \tilde{T}=\lim _{k \rightarrow+\infty} T \wedge d^{\prime}\left(\theta \circ v_{k}\right)
$$

in the weak topology of currents. It is therefore sufficient to check that $T \wedge d^{\prime}\left(\theta \circ v_{k}\right)$ converges weakly to 0 in $\mathscr{D}_{p-1, p}^{\prime}(\Omega)$ (note that $d^{\prime \prime} \tilde{T}$ is conjugate to $d^{\prime} \tilde{T}$, thus $d^{\prime \prime} \tilde{T}$ will also vanish).

Assume first that $p=1$. Then $T \wedge d^{\prime}\left(\theta \circ v_{k}\right) \in \mathscr{D}_{0,1}^{\prime}(\Omega)$, and we have to show that

$$
\left\langle T \wedge d^{\prime}\left(\theta \circ v_{k}\right), \bar{\alpha}\right\rangle=\left\langle T, \theta^{\prime}\left(v_{k}\right) d^{\prime} v_{k} \wedge \bar{\alpha}\right\rangle \longrightarrow 0, \quad \forall \alpha \in \mathscr{D}_{1,0}(\Omega)
$$

As $\gamma \longmapsto\langle T, \mathrm{i} \gamma \wedge \bar{\gamma}\rangle$ is a non-negative hermitian form on $\mathscr{D}_{1,0}(\Omega)$, the Cauchy-Schwarz inequality yields

$$
|\langle T, \mathrm{i} \beta \wedge \bar{\gamma}\rangle|^{2} \leqslant\langle T, \mathrm{i} \beta \wedge \bar{\beta}\rangle\langle T, \mathrm{i} \gamma \wedge \bar{\gamma}\rangle, \quad \forall \beta, \gamma \in \mathscr{D}_{1,0}(\Omega)
$$

Let $\psi \in \mathscr{D}(\Omega), 0 \leqslant \psi \leqslant 1$, be equal to 1 in a neighborhood of $\operatorname{Supp} \alpha$. We find

$$
\left|\left\langle T, \theta^{\prime}\left(v_{k}\right) d^{\prime} v_{k} \wedge \bar{\alpha}\right\rangle\right|^{2} \leqslant\left\langle T, \psi \mathrm{i} d^{\prime} v_{k} \wedge d^{\prime \prime} v_{k}\right\rangle\left\langle T, \theta^{\prime}\left(v_{k}\right)^{2} \mathrm{i} \alpha \wedge \bar{\alpha}\right\rangle
$$

By hypothesis $\int_{\Omega \backslash E} T \wedge \mathrm{i} \alpha \wedge \bar{\alpha}<+\infty$ and $\theta^{\prime}\left(v_{k}\right)$ converges everywhere to 0 on $\Omega$, thus $\left\langle T, \theta^{\prime}\left(v_{k}\right)^{2} \mathrm{i} \alpha \wedge \bar{\alpha}\right\rangle$ converges to 0 by Lebesgue's dominated convergence theorem. On the other hand

$$
\begin{aligned}
& \mathrm{i} d^{\prime} d^{\prime \prime} v_{k}^{2}=2 v_{k} \mathrm{i} d^{\prime} d^{\prime \prime} v_{k}+2 \mathrm{i} d^{\prime} v_{k} \wedge d^{\prime \prime} v_{k} \geqslant 2 \mathrm{i} d^{\prime} v_{k} \wedge d^{\prime \prime} v_{k} \\
& 2\left\langle T, \psi \mathrm{i} d^{\prime} v_{k} \wedge d^{\prime \prime} v_{k}\right\rangle \leqslant\left\langle T, \psi \mathrm{i} d^{\prime} d^{\prime \prime} v_{k}^{2}\right\rangle
\end{aligned}
$$

As $\psi \in \mathscr{D}(\Omega), v_{k}=0$ near $E$ and $d^{\prime} T=d^{\prime \prime} T=0$ on $\Omega \backslash E$, an integration by parts yields

$$
\left\langle T, \psi \mathrm{i} d^{\prime} d^{\prime \prime} v_{k}^{2}\right\rangle=\left\langle T, v_{k}^{2} \mathrm{i} d^{\prime} d^{\prime \prime} \psi\right\rangle \leqslant C \int_{\Omega \backslash E}\|T\|<+\infty
$$

where $C$ is a bound for the coefficients of $\psi$. Thus $\left\langle T, \psi \mathrm{i} d^{\prime} v_{k} \wedge d^{\prime \prime} v_{k}\right\rangle$ is bounded, and the proof is complete when $p=1$.

In the general case, let $\beta_{s}=\mathrm{i} \beta_{s, 1} \wedge \bar{\beta}_{s, 1} \wedge \ldots \wedge \mathrm{i} \beta_{s, p-1} \wedge \bar{\beta}_{s, p-1}$ be a basis of forms of bidegree ( $p-1, p-1$ ) with constant coefficients (Lemma 1.4). Then $T \wedge \beta_{s} \in \mathscr{D}_{1,1}^{\prime+}(\Omega \backslash E)$
has finite mass near $E$ and is closed on $\Omega \backslash E$. Therefore $d\left(\tilde{T} \wedge \beta_{s}\right)=(d \tilde{T}) \wedge \beta_{s}=0$ on $\Omega$ for all $s$, and we conclude that $d \tilde{T}=0$.
(2.4) Corollary. If $T \in \mathscr{D}_{p, p}^{\prime+}(X)$ is closed, if $E \subset X$ is a closed complete pluripolar set and $\mathbb{1}_{E}$ is its characteristic function, then $\mathbb{1}_{E} T$ and $\mathbb{1}_{X \backslash E} T$ are closed (and, of course, positive).

Proof. If we set $\Theta=T_{\upharpoonright X \backslash E}$, then $\Theta$ has finite mass near $E$ and we have $\mathbb{1}_{X \backslash E} T=\tilde{\Theta}$ and $\mathbb{1}_{E} T=T-\tilde{\Theta}$.

## § 2.B. Current of Integration over an Analytic Set

Let $A$ be a pure $p$-dimensional analytic subset of a complex manifold $X$. We would like to generalize Example 1.20 and to define a current of integration $[A]$ by letting

$$
\begin{equation*}
\langle[A], v\rangle=\int_{A_{\mathrm{reg}}} v, \quad v \in \mathscr{D}_{p, p}(X) \tag{2.5}
\end{equation*}
$$

One difficulty is of course to verify that the integral converges near $A_{\text {sing }}$. This follows from the following lemma, due to [Lelong 1957].
(2.6) Lemma. The current $\left[A_{\text {reg }}\right] \in \mathscr{D}_{p, p}^{++}\left(X \backslash A_{\text {sing }}\right)$ has finite mass in a neighborhood of every point $z_{0} \in A_{\text {sing }}$.

Proof. Set $T=\left[A_{\text {reg }}\right]$ and let $\Omega \ni z_{0}$ be a coordinate open set. If we write the monomials $d z_{K} \wedge d \bar{z}_{L}$ in terms of an arbitrary basis of $\Lambda^{p, p} T^{\star} \Omega$ consisting of decomposable forms $\beta_{s}=\mathrm{i} \beta_{s, 1} \wedge \bar{\beta}_{s, 1} \wedge \ldots \wedge \beta_{s, p} \wedge \bar{\beta}_{s, p}$ (cf. Lemma 1.4), we see that the measures $T_{I, J} . \tau$ are linear combinations of the positive measures $T \wedge \beta_{s}$. It is thus sufficient to prove that all $T \wedge \beta_{s}$ have finite mass near $A_{\text {sing }}$. Without loss of generality, we may assume that $\left(A, z_{0}\right)$ is irreducible. Take new coordinates $w=\left(w_{1}, \ldots, w_{n}\right)$ such that $w_{j}=\beta_{s, j}\left(z-z_{0}\right)$, $1 \leqslant j \leqslant p$. After a slight perturbation of the $\beta_{s, j}$, we may assume that each projection

$$
\pi_{s}: A \cap\left(\Delta^{\prime} \times \Delta^{\prime \prime}\right), \quad w \longmapsto w^{\prime}=\left(w_{1}, \ldots, w_{p}\right)
$$

defines a ramified covering of $A$ (cf. Prop. II.3.8 and Th. II.3.19), and that ( $\beta_{s}$ ) remains a basis of $\Lambda^{p, p} T^{\star} \Omega$. Let $S$ be the ramification locus of $\pi_{s}$ and $A_{S}=A \cap\left(\left(\Delta^{\prime} \backslash S\right) \times \Delta^{\prime \prime}\right) \subset$ $A_{\text {reg. }}$. The restriction of $\pi_{s}: A_{S} \longrightarrow \Delta^{\prime} \backslash S$ is then a covering with finite sheet number $q_{s}$ and we find

$$
\begin{aligned}
& \int_{\Delta^{\prime} \times \Delta^{\prime \prime}}\left[A_{\mathrm{reg}}\right] \wedge \beta_{s}=\int_{A_{\mathrm{reg} \cap\left(\Delta^{\prime} \times \Delta^{\prime \prime}\right)} \mathrm{i} d w_{1} \wedge d \bar{w}_{1} \wedge \ldots \wedge \mathrm{i} d w_{p} \wedge d \bar{w}_{p}}^{\quad=\int_{A_{S}} \mathrm{i} d w_{1} \wedge d \bar{w}_{1} \ldots \wedge d \bar{w}_{p}=q_{s} \int_{\Delta^{\prime} \backslash S} \mathrm{i} d w_{1} \wedge d \bar{w}_{1} \ldots \wedge d \bar{w}_{p}<+\infty} .
\end{aligned}
$$

The second equality holds because $A_{S}$ is the complement in $A_{\text {reg }} \cap\left(\Delta^{\prime} \times \Delta^{\prime \prime}\right)$ of an analytic subset (such a set is of zero Lebesgue measure in $A_{\text {reg }}$ ).
(2.7) Theorem ([Lelong 1957]). For every pure p-dimensional analytic subset $A \subset X$, the current of integration $[A] \in \mathscr{D}_{p, p}^{\prime+}(X)$ is a closed positive current on $X$.

Proof. Indeed, $\left[A_{\text {reg }}\right]$ has finite mass near $A_{\text {sing }}$ and $[A]$ is the trivial extension of $\left[A_{\text {reg }}\right]$ to $X$ through the complete pluripolar set $E=A_{\text {sing }}$. Theorem 2.7 is then a consequence of El Mir's theorem.

## § 2.C. Support Theorems and Lelong-Poincaré Equation

Let $M \subset X$ be a closed $C^{1}$ real submanifold of $X$. The holomorphic tangent space at a point $x \in M$ is

$$
\begin{equation*}
{ }^{h} T_{x} M=T_{x} M \cap J T_{x} M, \tag{2.8}
\end{equation*}
$$

that is, the largest complex subspace of $T_{x} X$ contained in $T_{x} M$. We define the CauchyRiemann dimension of $M$ at $x$ by $\operatorname{CRdim}_{x} M=\operatorname{dim}_{\mathbb{C}}{ }^{h} T_{x} M$ and say that $M$ is a $C R$ submanifold of $X$ if $\mathrm{CR}_{\operatorname{dim}_{x}} M$ is a constant. In general, we set

$$
\begin{equation*}
\text { CRdim } M=\max _{x \in M} \operatorname{CRdim}_{x} M=\max _{x \in M} \operatorname{dim}_{\mathbb{C}}^{h} T_{x} M \tag{2.9}
\end{equation*}
$$

A current $\Theta$ is said to be normal if $\Theta$ and $d \Theta$ are currents of order 0 . For instance, every closed positive current is normal. We are going to prove two important theorems describing the structure of normal currents with support in $C R$ submanifolds.
(2.10) First theorem of support. Let $\Theta \in \mathscr{D}_{p, p}^{\prime}(X)$ be a normal current. If Supp $\Theta$ is contained in a real submanifold $M$ of $C R$ dimension $<p$, then $\Theta=0$.

Proof. Let $x_{0} \in M$ and let $g_{1}, \ldots, g_{m}$ be real $C^{1}$ functions in a neighborhood $\Omega$ of $x_{0}$ such that $M=\left\{z \in \Omega ; g_{1}(z)=\ldots=g_{m}(z)=0\right\}$ and $d g_{1} \wedge \ldots \wedge d g_{m} \neq 0$ on $\Omega$. Then

$$
{ }^{h} T M=T M \cap J T M=\bigcap_{1 \leqslant k \leqslant m} \operatorname{ker} d g_{k} \cap \operatorname{ker}\left(d g_{k} \circ J\right)=\bigcap_{1 \leqslant k \leqslant m} \operatorname{ker} d^{\prime} g_{k}
$$

because $d^{\prime} g_{k}=\frac{1}{2}\left(d g_{k}-i\left(d g_{k}\right) \circ J\right)$. As $\operatorname{dim}_{\mathbb{C}}{ }^{h} T M<p$, the rank of the system of $(1,0)$ forms $\left(d^{\prime} g_{k}\right)$ must be $>n-p$ at every point of $M \cap \Omega$. After a change of the ordering, we may assume for example that $\zeta_{1}=d^{\prime} g_{1}, \zeta_{2}=d^{\prime} g_{2}, \ldots, \zeta_{n-p+1}=d^{\prime} g_{n-p+1}$ are linearly independent on $\Omega$ (shrink $\Omega$ if necessary). Complete ( $\zeta_{1}, \ldots, \zeta_{n-p+1}$ ) into a continuous frame $\left(\zeta_{1}, \ldots, \zeta_{n}\right)$ of $T^{\star} X_{\mid \Omega}$ and set

$$
\Theta=\sum_{|I|=|J|=n-p} \Theta_{I, J} \zeta_{I} \wedge \bar{\zeta}_{J} \quad \text { on } \Omega
$$

As $\Theta$ and $d^{\prime} \Theta$ have measure coefficients supported on $M$ and $g_{k}=0$ on $M$, we get $g_{k} \Theta=g_{k} d^{\prime} \Theta=0$, thus

$$
d^{\prime} g_{k} \wedge \Theta=d^{\prime}\left(g_{k} \Theta\right)-g_{k} d^{\prime} \Theta=0, \quad 1 \leqslant k \leqslant m
$$

in particular $\zeta_{k} \wedge \Theta=0$ for all $1 \leqslant k \leqslant n-p+1$. When $|I|=n-p$, the multi-index $\complement I$ contains at least one of the elements $1, \ldots, n-p+1$, hence $\Theta \wedge \zeta_{C_{I}} \wedge \bar{\zeta}_{C J}=0$ and $\Theta_{I, J}=0$.
(2.11) Corollary. Let $\Theta \in \mathscr{D}_{p, p}^{\prime}(X)$ be a normal current. If Supp $\Theta$ is contained in an analytic subset $A$ of dimension $<p$, then $\Theta=0$.

Proof. As $A_{\text {reg }}$ is a submanifold of CRdim $<p$ in $X \backslash A_{\text {sing }}$, Theorem 2.9 shows that Supp $\Theta \subset A_{\text {sing }}$ and we conclude by induction on $\operatorname{dim} A$.

Now, assume that $M \subset X$ is a CR submanifold of class $C^{1}$ with CRdim $M=p$ and that ${ }^{h} T M$ is an integrable subbundle of $T M$; this means that the Lie bracket of two vector fields in ${ }^{h} T M$ is in ${ }^{h} T M$. The Frobenius integrability theorem then shows that $M$ is locally fibered by complex analytic $p$-dimensional submanifolds. More precisely, in a neighborhood of every point of $M$, there is a submersion $\sigma: M \longrightarrow Y$ onto a real $C^{1}$ manifold $Y$ such that the tangent space to each fiber $F_{t}=\sigma^{-1}(t), t \in Y$, is the holomorphic tangent space ${ }^{h} T M$; moreover, the fibers $F_{t}$ are necessarily complex analytic in view of Lemma 1.7.18. Under these assumptions, with any complex measure $\mu$ on $Y$ we associate a current $\Theta$ with support in $M$ by

$$
\begin{equation*}
\Theta=\int_{t \in Y}\left[F_{t}\right] d \mu(t), \quad \text { i.e. }\langle\Theta, u\rangle=\int_{t \in Y}\left(\int_{F_{t}} u\right) d \mu(t) \tag{2.12}
\end{equation*}
$$

for all $u \in \mathscr{D}_{p, p}^{\prime}(X)$. Then clearly $\Theta \in \mathscr{D}_{p, p}^{\prime}(X)$ is a closed current of order 0 , for all fibers $\left[F_{t}\right]$ have the same properties. When the fibers $F_{t}$ are connected, the following converse statement holds:
(2.13) Second theorem of support. Let $M \subset X$ be a $C R$ submanifold of $C R$ dimension $p$ such that there is a submersion $\sigma: M \longrightarrow Y$ of class $C^{1}$ whose fibers $F_{t}=\sigma^{-1}(t)$ are connected and are the integral manifolds of the holomorphic tangent space ${ }^{h} T M$. Then any closed current $\Theta \in \mathscr{D}_{p, p}^{\prime}(X)$ of order 0 with support in $M$ can be written $\Theta=\int_{Y}\left[F_{t}\right] d \mu(t)$ with a unique complex measure $\mu$ on $Y$. Moreover $\Theta$ is (strongly) positive if and only if the measure $\mu$ is positive.

Proof. Fix a compact set $K \subset Y$ and a $C^{1}$ retraction $\rho$ from a neighborhood $V$ of $M$ onto $M$. By means of a partition of unity, it is easy to construct a positive form $\alpha \in \mathscr{D}_{p, p}^{0}(V)$ such that $\int_{F_{t}} \alpha=1$ for each fiber $F_{t}$ with $t \in K$. Then the uniqueness and positivity statements for $\mu$ follow from the obvious formula

$$
\int_{Y} f(t) d \mu(t)=\langle\Theta,(f \circ \sigma \circ \rho) \alpha\rangle, \quad \forall f \in C^{0}(Y), \quad \operatorname{Supp} f \subset K
$$

Now, let us prove the existence of $\mu$. Let $x_{0} \in M$. There is a small neighborhood $\Omega$ of $x_{0}$ and real coordinates $\left(x_{1}, y_{1}, \ldots, x_{p}, y_{p}, t_{1}, \ldots, t_{q}, g_{1}, \ldots, g_{m}\right)$ such that

- $z_{j}=x_{j}+\mathrm{i} y_{j}, 1 \leqslant j \leqslant p$, are holomorphic functions on $\Omega$ that define complex coordinates on all fibers $F_{t} \cap \Omega$.
- $t_{1}, \ldots, t_{q}$ restricted to $M \cap \Omega$ are pull-backs by $\sigma: M \rightarrow Y$ of local coordinates on an open set $U \subset Y$ such that $\sigma_{\mid \Omega}: M \cap \Omega \longrightarrow U$ is a trivial fiber space.
- $g_{1}=\ldots=g_{m}=0$ are equations of $M$ in $\Omega$.

Then $T F_{t}=\left\{d t_{j}=d g_{k}=0\right\}$ equals ${ }^{h} T M=\left\{d^{\prime} g_{k}=0\right\}$ and the rank of $\left(d^{\prime} g_{1}, \ldots, d^{\prime} g_{m}\right)$ is equal to $n-p$ at every point of $M \cap \Omega$. After a change of the ordering we may suppose that $\zeta_{1}=d^{\prime} g_{1}, \ldots, \zeta_{n-p}=d^{\prime} g_{n-p}$ are linearly independent on $\Omega$. As in Prop. 2.10, we get $\zeta_{k} \wedge \Theta=\bar{\zeta}_{k} \wedge \Theta=0$ for $1 \leqslant k \leqslant n-p$ and infer that $\Theta \wedge \zeta_{C_{I}} \wedge \bar{\zeta}_{C_{J}}=0$ unless $I=J=L$ where $L=\{1,2, \ldots, n-p\}$. Hence

$$
\Theta=\Theta_{L, L} \zeta_{1} \wedge \ldots \wedge \zeta_{n-p} \wedge \bar{\zeta}_{1} \wedge \ldots \wedge \bar{\zeta}_{n-p} \quad \text { on } \Omega
$$

Now $\zeta_{1} \wedge \ldots \wedge \bar{\zeta}_{n-p}$ is proportional to $d t_{1} \wedge \ldots d t_{q} \wedge d g_{1} \wedge \ldots \wedge d g_{m}$ because both induce a volume form on the quotient space $T X_{\lceil M} /{ }^{h} T M$. Therefore, there is a complex measure $\nu$ supported on $M \cap \Omega$ such that

$$
\Theta=\nu d t_{1} \wedge \ldots d t_{q} \wedge d g_{1} \wedge \ldots \wedge d g_{m} \quad \text { on } \Omega
$$

As $\Theta$ is supposed to be closed, we have $\partial \nu / \partial x_{j}=\partial \nu / \partial y_{j}=0$. Hence $\nu$ is a measure depending only on $(t, g)$, with support in $g=0$. We may write $\nu=d \mu_{U}(t) \otimes \delta_{0}(g)$ where $\mu_{U}$ is a measure on $U=\sigma(M \cap \Omega)$ and $\delta_{0}$ is the Dirac measure at 0 . If $j: M \longrightarrow X$ is the injection, this means precisely that $\Theta=j_{\star} \sigma^{\star} \mu_{U}$ on $\Omega$, i.e.

$$
\Theta=\int_{t \in U}\left[F_{t}\right] d \mu_{U}(t) \quad \text { on } \Omega
$$

The uniqueness statement shows that for two open sets $\Omega_{1}, \Omega_{2}$ as above, the associated measures $\mu_{U_{1}}$ and $\mu_{U_{2}}$ coincide on $\sigma\left(M \cap \Omega_{1} \cap \Omega_{2}\right)$. Since the fibers $F_{t}$ are connected, there is a unique measure $\mu$ which coincides with all measures $\mu_{U}$.
(2.14) Corollary. Let $A$ be an analytic subset of $X$ with global irreducible components $A_{j}$ of pure dimension $p$. Then any closed current $\Theta \in \mathscr{D}_{p, p}^{\prime}(X)$ of order 0 with support in $A$ is of the form $\Theta=\sum \lambda_{j}\left[A_{j}\right]$ where $\lambda_{j} \in \mathbb{C}$. Moreover, $\Theta$ is (strongly) positive if and only if all coefficients $\lambda_{j}$ are $\geqslant 0$.

Proof. The regular part $M=A_{\text {reg }}$ is a complex submanifold of $X \backslash A_{\text {sing }}$ and its connected components are $A_{j} \cap A_{\text {reg }}$. Thus, we may apply Th. 2.13 in the case where $Y$ is discrete to see that $\Theta=\sum \lambda_{j}\left[A_{j}\right]$ on $X \backslash A_{\text {sing }}$. Now $\operatorname{dim} A_{\text {sing }}<p$ and the difference $\Theta-\sum \lambda_{j}\left[A_{j}\right] \in \mathscr{D}_{p, p}^{\prime}(X)$ is a closed current of order 0 with support in $A_{\text {sing }}$, so this current must vanish by Cor. 2.11.
(2.15) Lelong-Poincaré equation. Let $f \in \mathscr{M}(X)$ be a meromorphic function which does not vanish identically on any connected component of $X$ and let $\sum m_{j} Z_{j}$ be the divisor of $f$. Then the function $\log |f|$ is locally integrable on $X$ and satisfies the equation

$$
\frac{\mathrm{i}}{\pi} d^{\prime} d^{\prime \prime} \log |f|=\sum m_{j}\left[Z_{j}\right]
$$

in the space $\mathscr{D}_{n-1, n-1}^{\prime}(X)$ of currents of bidimension $(n-1, n-1)$.
We refer to Sect. 2.6 for the definition of divisors, and especially to (2.6.14). Observe that if $f$ is holomorphic, then $\log |f| \in \operatorname{Psh}(X)$, the coefficients $m_{j}$ are positive integers and the right hand side is a positive current in $\mathscr{D}_{n-1, n-1}^{\prime+}(X)$.
Proof. Let $Z=\bigcup Z_{j}$ be the support of $\operatorname{div}(f)$. Observe that the sum in the right hand side is locally finite and that $d^{\prime} d^{\prime \prime} \log |f|$ is supported on $Z$, since

$$
d^{\prime} \log |f|^{2}=d^{\prime} \log (f \bar{f})=\frac{\bar{f} d f}{f \bar{f}}=\frac{d f}{f} \quad \text { on } \quad X \backslash Z .
$$

In a neighborhood $\Omega$ of a point $a \in Z_{j} \cap Z_{\text {reg }}$, we can find local coordinates $\left(w_{1}, \ldots, w_{n}\right)$ such that $Z_{j} \cap \Omega$ is given by the equation $w_{1}=0$. Then Th. 2.6.6 shows that $f$ can
be written $f(w)=u(w) w_{1}^{m_{j}}$ with an invertible holomorphic function $u$ on a smaller neighborhood $\Omega^{\prime} \subset \Omega$. Then we have

$$
\mathrm{i} d^{\prime} d^{\prime \prime} \log |f|=\mathrm{i} d^{\prime} d^{\prime \prime}\left(\log |u|+m_{j} \log \left|w_{1}\right|\right)=m_{j} \mathrm{i} d^{\prime} d^{\prime \prime} \log \left|w_{1}\right| .
$$

For $z \in \mathbb{C}$, Cor. I.3.4 implies

$$
\mathrm{i} d^{\prime} d^{\prime \prime} \log |z|^{2}=-\mathrm{i} d^{\prime \prime}\left(\frac{d z}{z}\right)=-\mathrm{i} \pi \delta_{0} d \bar{z} \wedge d z=2 \pi[0] .
$$

If $\Phi: \mathbb{C}^{n} \longrightarrow \mathbb{C}$ is the projection $z \longmapsto z_{1}$ and $H \subset \mathbb{C}^{n}$ the hyperplane $\left\{z_{1}=0\right\}$, formula (1.2.19) shows that

$$
\mathrm{i} d^{\prime} d^{\prime \prime} \log \left|z_{1}\right|=\mathrm{i} d^{\prime} d^{\prime \prime} \log |\Phi(z)|=\Phi^{\star}\left(\mathrm{i} d^{\prime} d^{\prime \prime} \log |z|\right)=\pi \Phi^{\star}([0])=\pi[H],
$$

because $\Phi$ is a submersion. We get therefore $\frac{\mathrm{i}}{\pi} d^{\prime} d^{\prime \prime} \log |f|=m_{j}\left[Z_{j}\right]$ in $\Omega^{\prime}$. This implies that the Lelong-Poincaré equation is valid at least on $X \backslash Z_{\operatorname{sing}}$. As $\operatorname{dim} Z_{\text {sing }}<n-1$, Cor. 2.11 shows that the equation holds everywhere on $X$.

## § 3. Definition of Monge-Ampère Operators

Let $X$ be a $n$-dimensional complex manifold. We denote by $d=d^{\prime}+d^{\prime \prime}$ the usual decomposition of the exterior derivative in terms of its $(1,0)$ and $(0,1)$ parts, and we set

$$
d^{c}=\frac{1}{2 \mathrm{i} \pi}\left(d^{\prime}-d^{\prime \prime}\right) .
$$

It follows in particular that $d^{c}$ is a real operator, i.e. $\overline{d^{c} u}=d^{c} \bar{u}$, and that $d d^{c}=\frac{\mathrm{i}}{\pi} d^{\prime} d^{\prime \prime}$. Although not quite standard, the $1 / 2 \mathrm{i} \pi$ normalization is very convenient for many purposes, since we may then forget the factor $2 \pi$ almost everywhere (e.g. in the Lelong-Poincaré equation (2.15)). In this context, we have the following integration by part formula.
(3.1) Formula. Let $\Omega \subset \subset X$ be a smoothly bounded open set in $X$ and let $f, g$ be forms of class $C^{2}$ on $\bar{\Omega}$ of pure bidegrees $(p, p)$ and $(q, q)$ with $p+q=n-1$. Then

$$
\int_{\Omega} f \wedge d d^{c} g-d d^{c} f \wedge g=\int_{\partial \Omega} f \wedge d^{c} g-d^{c} f \wedge g
$$

Proof. By Stokes' theorem the right hand side is the integral over $\Omega$ of

$$
d\left(f \wedge d^{c} g-d^{c} f \wedge g\right)=f \wedge d d^{c} g-d d^{c} f \wedge g+\left(d f \wedge d^{c} g+d^{c} f \wedge d g\right)
$$

As all forms of total degree $2 n$ and bidegree $\neq(n, n)$ are zero, we get

$$
d f \wedge d^{c} g=\frac{1}{2 \mathrm{i} \pi}\left(d^{\prime \prime} f \wedge d^{\prime} g-d^{\prime} f \wedge d^{\prime \prime} g\right)=-d^{c} f \wedge d g
$$

Let $u$ be a plurisubharmonic function on $X$ and let $T$ be a closed positive current of bidimension $(p, p)$, i.e. of bidegree $(n-p, n-p)$. Our desire is to define the wedge
product $d d^{c} u \wedge T$ even when neither $u$ nor $T$ are smooth. A priori, this product does not make sense because $d d^{c} u$ and $T$ have measure coefficients and measures cannot be multiplied; see [Kiselman 1983] for interesting counterexamples. Assume however that $u$ is a locally bounded plurisubharmonic function. Then the current $u T$ is well defined since $u$ is a locally bounded Borel function and $T$ has measure coefficients. According to [Bedford-Taylor 1982] we define

$$
d d^{c} u \wedge T=d d^{c}(u T)
$$

where $d d^{c}()$ is taken in the sense of distribution (or current) theory.
(3.2) Proposition. The wedge product $d d^{c} u \wedge T$ is again a closed positive current.

Proof. The result is local. In an open set $\Omega \subset \mathbb{C}^{n}$, we can use convolution with a family of regularizing kernels to find a decreasing sequence of smooth plurisubharmonic functions $u_{k}=u \star \rho_{1 / k}$ converging pointwise to $u$. Then $u \leqslant u_{k} \leqslant u_{1}$ and Lebesgue's dominated convergence theorem shows that $u_{k} T$ converges weakly to $u T$; thus $d d^{c}\left(u_{k} T\right)$ converges weakly to $d d^{c}(u T)$ by the weak continuity of differentiations. However, since $u_{k}$ is smooth, $d d^{c}\left(u_{k} T\right)$ coincides with the product $d d^{c} u_{k} \wedge T$ in its usual sense. As $T \geqslant 0$ and as $d d^{c} u_{k}$ is a positive (1,1)-form, we have $d d^{c} u_{k} \wedge T \geqslant 0$, hence the weak limit $d d^{c} u \wedge T$ is $\geqslant 0$ (and obviously closed).

Given locally bounded plurisubharmonic functions $u_{1}, \ldots, u_{q}$, we define inductively

$$
d d^{c} u_{1} \wedge d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T=d d^{c}\left(u_{1} d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T\right)
$$

By (3.2) the product is a closed positive current. In particular, when $u$ is a locally bounded plurisubharmonic function, the bidegree $(n, n)$ current $\left(d d^{c} u\right)^{n}$ is well defined and is a positive measure. If $u$ is of class $C^{2}$, a computation in local coordinates gives

$$
\left(d d^{c} u\right)^{n}=\operatorname{det}\left(\frac{\partial^{2} u}{\partial z_{j} \partial \bar{z}_{k}}\right) \cdot \frac{n!}{\pi^{n}} \mathrm{i} d z_{1} \wedge d \bar{z}_{1} \wedge \ldots \wedge \mathrm{i} d z_{n} \wedge d \bar{z}_{n}
$$

The expression "Monge-Ampère operator" classically refers to the non-linear partial differential operator $u \longmapsto \operatorname{det}\left(\partial^{2} u / \partial z_{j} \partial \bar{z}_{k}\right)$. By extension, all operators $\left(d d^{c}\right)^{q}$ defined above are also called Monge-Ampère operators.

Now, let $\Theta$ be a current of order 0 . When $K \subset \subset X$ is an arbitrary compact subset, we define a mass semi-norm

$$
\|\Theta\|_{K}=\sum_{j} \int_{K_{j}} \sum_{I, J}\left|\Theta_{I, J}\right|
$$

by taking a partition $K=\bigcup K_{j}$ where each $\bar{K}_{j}$ is contained in a coordinate patch and where $\Theta_{I, J}$ are the corresponding measure coefficients. Up to constants, the semi-norm $\|\Theta\|_{K}$ does not depend on the choice of the coordinate systems involved. When $K$ itself is contained in a coordinate patch, we set $\beta=d d^{c}|z|^{2}$ over $K$; then, if $\Theta \geqslant 0$, there are constants $C_{1}, C_{2}>0$ such that

$$
C_{1}\|\Theta\|_{K} \leqslant \int_{K} \Theta \wedge \beta^{p} \leqslant C_{2}\|\Theta\|_{K}
$$

We denote by $L^{1}(K)$, resp. by $L^{\infty}(K)$, the space of integrable (resp. bounded measurable) functions on $K$ with respect to any smooth positive density on $X$.
(3.3) Chern-Levine-Nirenberg inequalities (1969). For all compact subsets $K, L$ of $X$ with $L \subset K^{\circ}$, there exists a constant $C_{K, L} \geqslant 0$ such that

$$
\left\|d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T\right\|_{L} \leqslant C_{K, L}\left\|u_{1}\right\|_{L^{\infty}(K)} \ldots\left\|u_{q}\right\|_{L^{\infty}(K)}\|T\|_{K}
$$

Proof. By induction, it is sufficient to prove the result for $q=1$ and $u_{1}=u$. There is a covering of $L$ by a family of balls $B_{j}^{\prime} \subset \subset B_{j} \subset K$ contained in coordinate patches of $X$. Let $\chi \in \mathscr{D}\left(B_{j}\right)$ be equal to 1 on $\bar{B}_{j}^{\prime}$. Then

$$
\left\|d d^{c} u \wedge T\right\|_{L \cap \bar{B}_{j}^{\prime}} \leqslant C \int_{\bar{B}_{j}^{\prime}} d d^{c} u \wedge T \wedge \beta^{p-1} \leqslant C \int_{B_{j}} \chi d d^{c} u \wedge T \wedge \beta^{p-1}
$$

As $T$ and $\beta$ are closed, an integration by parts yields

$$
\left\|d d^{c} u \wedge T\right\|_{L \cap \bar{B}_{j}^{\prime}} \leqslant C \int_{B_{j}} u T \wedge d d^{c} \chi \wedge \beta^{p-1} \leqslant C^{\prime}\|u\|_{L^{\infty}(K)}\|T\|_{K}
$$

where $C^{\prime}$ is equal to $C$ multiplied by a bound for the coefficients of the smooth form $d d^{c} \chi \wedge \beta^{p-1}$.
(3.4) Remark. With the same notations as above, any plurisubharmonic function $V$ on $X$ satisfies inequalities of the type
a) $\left\|d d^{c} V\right\|_{L} \leqslant C_{K, L}\|V\|_{L^{1}(K)}$.
b) $\sup _{L} V_{+} \leqslant C_{K, L}\|V\|_{L^{1}(K)}$.

In fact the inequality

$$
\int_{L \cap \bar{B}_{j}^{\prime}} d d^{c} V \wedge \beta^{n-1} \leqslant \int_{B_{j}} \chi d d^{c} V \wedge \beta^{n-1}=\int_{B_{j}} V d d^{c} \chi \wedge \beta^{n-1}
$$

implies a), and b) follows from the mean value inequality.
(3.5) Remark. Products of the form $\Theta=\gamma_{1} \wedge \ldots \wedge \gamma_{q} \wedge T$ with mixed (1,1)-forms $\gamma_{j}=d d^{c} u_{j}$ or $\gamma_{j}=d v_{j} \wedge d^{c} w_{j}+d w_{j} \wedge d^{c} v_{j}$ are also well defined whenever $u_{j}, v_{j}, w_{j}$ are locally bounded plurisubharmonic functions. Moreover, for $L \subset K^{\circ}$, we have

$$
\|\Theta\|_{L} \leqslant C_{K, L}\|T\|_{K} \prod\left\|u_{j}\right\|_{L^{\infty}(K)} \prod\left\|v_{j}\right\|_{L^{\infty}(K)} \prod\left\|w_{j}\right\|_{L^{\infty}(K)}
$$

To check this, we may suppose $v_{j}, w_{j} \geqslant 0$ and $\left\|v_{j}\right\|=\left\|w_{j}\right\|=1$ in $L^{\infty}(K)$. Then the inequality follows from (3.3) by the polarization identity

$$
2\left(d v_{j} \wedge d^{c} w_{j}+d w_{j} \wedge d^{c} v_{j}\right)=d d^{c}\left(v_{j}+w_{j}\right)^{2}-d d^{c} v_{j}^{2}-d d^{c} w_{j}^{2}-v_{j} d d^{c} w_{j}-w_{j} d d^{c} v_{j}
$$

in which all $d d^{c}$ operators act on plurisubharmonic functions.
(3.6) Corollary. Let $u_{1}, \ldots, u_{q}$ be continuous (finite) plurisubharmonic functions and let $u_{1}^{k}, \ldots, u_{q}^{k}$ be sequences of plurisubharmonic functions converging locally uniformly to $u_{1}, \ldots, u_{q}$. If $T_{k}$ is a sequence of closed positive currents converging weakly to $T$, then
a) $u_{1}^{k} d d^{c} u_{2}^{k} \wedge \ldots \wedge d d^{c} u_{q}^{k} \wedge T_{k} \longrightarrow u_{1} d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \quad$ weakly.
b) $d d^{c} u_{1}^{k} \wedge \ldots \wedge d d^{c} u_{q}^{k} \wedge T_{k} \longrightarrow d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \quad$ weakly.

Proof. We observe that b) is an immediate consequence of a) by the weak continuity of $d d^{c}$. By using induction on $q$, it is enough to prove result a) when $q=1$. If ( $u^{k}$ ) converges locally uniformly to a finite continuous plurisubharmonic function $u$, we introduce local regularizations $u_{\varepsilon}=u \star \rho_{\varepsilon}$ and write

$$
u^{k} T_{k}-u T=\left(u^{k}-u\right) T_{k}+\left(u-u_{\varepsilon}\right) T_{k}+u_{\varepsilon}\left(T_{k}-T\right) .
$$

As the sequence $T_{k}$ is weakly convergent, it is locally uniformly bounded in mass, thus $\left\|\left(u^{k}-u\right) T_{k}\right\|_{K} \leqslant\left\|u^{k}-u\right\|_{L^{\infty}(K)}\left\|T_{k}\right\|_{K}$ converges to 0 on every compact set $K$. The same argument shows that $\left\|\left(u-u_{\varepsilon}\right) T_{k}\right\|_{K}$ can be made arbitrarily small by choosing $\varepsilon$ small enough. Finally $u_{\varepsilon}$ is smooth, so $u_{\varepsilon}\left(T_{k}-T\right)$ converges weakly to 0 .

Now, we prove a deeper monotone continuity theorem due to [Bedford-Taylor 1982] according to which the continuity and uniform convergence assumptions can be dropped if each sequence $\left(u_{j}^{k}\right)$ is decreasing and $T_{k}$ is a constant sequence.
(3.7) Theorem. Let $u_{1}, \ldots, u_{q}$ be locally bounded plurisubharmonic functions and let $u_{1}^{k}, \ldots, u_{q}^{k}$ be decreasing sequences of plurisubharmonic functions converging pointwise to $u_{1}, \ldots, u_{q}$. Then
a) $u_{1}^{k} d d^{c} u_{2}^{k} \wedge \ldots \wedge d d^{c} u_{q}^{k} \wedge T \longrightarrow u_{1} d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \quad$ weakly.
b) $d d^{c} u_{1}^{k} \wedge \ldots \wedge d d^{c} u_{q}^{k} \wedge T \longrightarrow d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \quad$ weakly.

Proof. Again by induction, observing that a$) \Longrightarrow \mathrm{b}$ ) and that a) is obvious for $q=1$ thanks to Lebesgue's bounded convergence theorem. To proceed with the induction step, we first have to make some slight modifications of our functions $u_{j}$ and $u_{j}^{k}$.

As the sequence $\left(u_{j}^{k}\right)$ is decreasing and as $u_{j}$ is locally bounded, the family $\left(u_{j}^{k}\right)_{k \in \mathbb{N}}$ is locally uniformly bounded. The results are local, so we can work on a Stein open set $\Omega \subset \subset X$ with strongly pseudoconvex boundary. We use the following notations:
let $\psi$ be a strongly plurisubharmonic function of class $\mathscr{C}^{\infty}$ near $\bar{\Omega}$ with $\psi<0$ on $\Omega$ and $\psi=0, d \psi \neq 0$ on $\partial \Omega$;

$$
\text { we set } \Omega_{\delta}=\{z \in \Omega ; \psi(z)<-\delta\} \text { for all } \delta>0 \text {. }
$$



III-1 Construction of $v_{j}^{k}$

After addition of a constant we can assume that $-M \leqslant u_{j}^{k} \leqslant-1$ near $\bar{\Omega}$. Let us denote by $\left.\left.\left(u_{j}^{k, \varepsilon}\right), \varepsilon \in\right] 0, \varepsilon_{0}\right]$, an increasing family of regularizations converging to $u_{j}^{k}$ as $\varepsilon \rightarrow 0$ and such that $-M \leqslant u_{j}^{k, \varepsilon} \leqslant-1$ on $\bar{\Omega}$. Set $A=M / \delta$ with $\delta>0$ small and replace $u_{j}^{k}$ by $v_{j}^{k}=\max \left\{A \psi, u_{j}^{k}\right\}$ and $u_{j}^{k, \varepsilon}$ by $v_{j}^{k, \varepsilon}=\max _{\varepsilon}\left\{A \psi, u_{j}^{k, \varepsilon}\right\}$ where $\max _{\varepsilon}=\max \star \rho_{\varepsilon}$ is a regularized max function.

Then $v_{j}^{k}$ coincides with $u_{j}^{k}$ on $\Omega_{\delta}$ since $A \psi<-A \delta=-M$ on $\Omega_{\delta}$, and $v_{j}^{k}$ is equal to $A \psi$ on the corona $\Omega \backslash \Omega_{\delta / M}$. Without loss of generality, we can therefore assume that all $u_{j}^{k}$ (and similarly all $u_{j}^{k, \varepsilon}$ ) coincide with $A \psi$ on a fixed neighborhood of $\partial \Omega$. We need a lemma.
(3.9) Lemma. Let $f_{k}$ be a decreasing sequence of upper semi-continuous functions converging to $f$ on some separable locally compact space $X$ and $\mu_{k}$ a sequence of positive measures converging weakly to $\mu$ on $X$. Then every weak limit $\nu$ of $f_{k} \mu_{k}$ satisfies $\nu \leqslant f \mu$.

Indeed if $\left(g_{p}\right)$ is a decreasing sequence of continuous functions converging to $f_{k_{0}}$ for some $k_{0}$, then $f_{k} \mu_{k} \leqslant f_{k_{0}} \mu_{k} \leqslant g_{p} \mu_{k}$ for $k \geqslant k_{0}$, thus $\nu \leqslant g_{p} \mu$ as $k \rightarrow+\infty$. The monotone convergence theorem then gives $\nu \leqslant f_{k_{0}} \mu$ as $p \rightarrow+\infty$ and $\nu \leqslant f \mu$ as $k_{0} \rightarrow+\infty$.

Proof of Theorem 3.7 (end).. Assume that a) has been proved for $q-1$. Then

$$
S^{k}=d d^{c} u_{2}^{k} \wedge \ldots \wedge d d^{c} u_{q}^{k} \wedge T \longrightarrow S=d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T
$$

By 3.3 the sequence $\left(u_{1}^{k} S^{k}\right)$ has locally bounded mass, hence is relatively compact for the weak topology. In order to prove a), we only have to show that every weak limit $\Theta$ of $u_{1}^{k} S^{k}$ is equal to $u_{1} S$. Let $(m, m)$ be the bidimension of $S$ and let $\gamma$ be an arbitrary smooth and strongly positive form of bidegree $(m, m)$. Then the positive measures $S^{k} \wedge \gamma$ converge weakly to $S \wedge \gamma$ and Lemma 3.9 shows that $\Theta \wedge \gamma \leqslant u_{1} S \wedge \gamma$, hence $\Theta \leqslant u_{1} S$. To get the equality, we set $\beta=d d^{c} \psi>0$ and show that $\int_{\Omega} u_{1} S \wedge \beta^{m} \leqslant \int_{\Omega} \Theta \wedge \beta^{m}$, i.e.

$$
\int_{\Omega} u_{1} d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \wedge \beta^{m} \leqslant \liminf _{k \rightarrow+\infty} \int_{\Omega} u_{1}^{k} d d^{c} u_{2}^{k} \wedge \ldots \wedge d d^{c} u_{q}^{k} \wedge T \wedge \beta^{m}
$$

As $u_{1} \leqslant u_{1}^{k} \leqslant u_{1}^{k, \varepsilon_{1}}$ for every $\varepsilon_{1}>0$, we get

$$
\begin{aligned}
& \int_{\Omega} u_{1} d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \wedge \beta^{m} \\
& \leqslant \int_{\Omega} u_{1}^{k, \varepsilon_{1}} d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \wedge \beta^{m} \\
&=\int_{\Omega} d d^{c} u_{1}^{k, \varepsilon_{1}} \wedge u_{2} d d^{c} u_{3} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \wedge \beta^{m}
\end{aligned}
$$

after an integration by parts (there is no boundary term because $u_{1}^{k, \varepsilon_{1}}$ and $u_{2}$ both vanish on $\partial \Omega$ ). Repeating this argument with $u_{2}, \ldots, u_{q}$, we obtain

$$
\begin{aligned}
& \int_{\Omega} u_{1} d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \wedge \beta^{m} \\
& \leqslant \int_{\Omega} d d^{c} u_{1}^{k, \varepsilon_{1}} \wedge \ldots \wedge d d^{c} u_{q-1}^{k, \varepsilon_{q-1}} \wedge u_{q} T \wedge \beta^{m} \\
& \leqslant \int_{\Omega} u_{1}^{k, \varepsilon_{1}} d d^{c} u_{2}^{k, \varepsilon_{2}} \wedge \ldots \wedge d d^{c} u_{q}^{k, \varepsilon_{q}} \wedge T \wedge \beta^{m}
\end{aligned}
$$

Now let $\varepsilon_{q} \rightarrow 0, \ldots, \varepsilon_{1} \rightarrow 0$ in this order. We have weak convergence at each step and $u_{1}^{k, \varepsilon_{1}}=0$ on the boundary; therefore the integral in the last line converges and we get the desired inequality

$$
\int_{\Omega} u_{1} d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \wedge \beta^{m} \leqslant \int_{\Omega} u_{1}^{k} d d^{c} u_{2}^{k} \wedge \ldots \wedge d d^{c} u_{q}^{k} \wedge T \wedge \beta^{m}
$$

(3.10) Corollary. The product $d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T$ is symmetric with respect to $u_{1}, \ldots, u_{q}$.

Proof. Observe that the definition was unsymmetric. The result is true when $u_{1}, \ldots, u_{q}$ are smooth and follows in general from Th. 3.7 applied to the sequences $u_{j}^{k}=u_{j} \star \rho_{1 / k}$, $1 \leqslant j \leqslant q$.
(3.11) Proposition. Let $K, L$ be compact subsets of $X$ such that $L \subset K^{\circ}$. For any plurisubharmonic functions $V, u_{1}, \ldots, u_{q}$ on $X$ such that $u_{1}, \ldots, u_{q}$ are locally bounded, there is an inequality

$$
\left\|V d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q}\right\|_{L} \leqslant C_{K, L}\|V\|_{L^{1}(K)}\left\|u_{1}\right\|_{L^{\infty}(K)} \ldots\left\|u_{q}\right\|_{L^{\infty}(K)}
$$

Proof. We may assume that $L$ is contained in a strongly pseudoconvex open set $\Omega=\{\psi<$ $0\} \subset K$ (otherwise we cover $L$ by small balls contained in $K$ ). A suitable normalization gives $-2 \leqslant u_{j} \leqslant-1$ on $K$; then we can modify $u_{j}$ on $\Omega \backslash L$ so that $u_{j}=A \psi$ on $\Omega \backslash \Omega_{\delta}$ with a fixed constant $A$ and $\delta>0$ such that $L \subset \Omega_{\delta}$. Let $\chi \geqslant 0$ be a smooth function equal to $-\psi$ on $\Omega_{\delta}$ with compact support in $\Omega$. If we take $\|V\|_{L^{1}(K)}=1$, we see that $V_{+}$is uniformly bounded on $\Omega_{\delta}$ by 3.4 b ); after subtraction of a fixed constant we can
assume $V \leqslant 0$ on $\Omega_{\delta}$. First suppose $q \leqslant n-1$. As $u_{j}=A \psi$ on $\Omega \backslash \Omega_{\delta}$, we find

$$
\begin{aligned}
& \int_{\Omega_{\delta}}-V d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge \beta^{n-q} \\
& =\int_{\Omega} V d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge \beta^{n-q-1} \wedge d d^{c} \chi-A^{q} \int_{\Omega \backslash \Omega_{\delta}} V \beta^{n-1} \wedge d d^{c} \chi \\
& =\int_{\Omega} \chi d d^{c} V \wedge d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge \beta^{n-q-1}-A^{q} \int_{\Omega \backslash \Omega_{\delta}} V \beta^{n-1} \wedge d d^{c} \chi .
\end{aligned}
$$

The first integral of the last line is uniformly bounded thanks to 3.3 and 3.4 a), and the second one is bounded by $\|V\|_{L^{1}(\Omega)} \leqslant$ constant. Inequality 3.11 follows for $q \leqslant n-1$. If $q=n$, we can work instead on $X \times \mathbb{C}$ and consider $V, u_{1}, \ldots, u_{q}$ as functions on $X \times \mathbb{C}$ independent of the extra variable in $\mathbb{C}$.

## § 4. Case of Unbounded Plurisubharmonic Functions

We would like to define $d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T$ also in some cases when $u_{1}, \ldots, u_{q}$ are not bounded below everywhere, especially when the $u_{j}$ have logarithmic poles. Consider first the case $q=1$ and let $u$ be a plurisubharmonic function on $X$. The pole set of $u$ is by definition $P(u)=u^{-1}(-\infty)$. We define the unbounded locus $L(u)$ to be the set of points $x \in X$ such that $u$ is unbounded in every neighborhood of $x$. Clearly $L(u)$ is closed and we have $L(u) \supset \overline{P(u)}$ but in general these sets are different: in fact, $u(z)=\sum k^{-2} \log \left(|z-1 / k|+e^{-k^{3}}\right)$ is everywhere finite in $\mathbb{C}$ but $L(u)=\{0\}$.
(4.1) Proposition. We make two additional assumptions:
a) $T$ has non zero bidimension $(p, p)$ (i.e. degree of $T<2 n)$.
b) $X$ is covered by a family of Stein open sets $\Omega \subset \subset X$ whose boundaries $\partial \Omega$ do not intersect $L(u) \cap$ Supp $T$.
Then the current uT has locally finite mass in $X$.

For any current $T$, hypothesis 4.1 b ) is clearly satisfied when $u$ has a discrete unbounded locus $L(u)$; an interesting example is $u=\log |F|$ where $F=\left(F_{1}, \ldots, F_{N}\right)$ are holomorphic functions having a discrete set of common zeros. Observe that the current $u T$ need not have locally finite mass when $T$ has degree $2 n$ (i.e. $T$ is a measure); example: $T=\delta_{0}$ and $u(z)=\log |z|$ in $\mathbb{C}^{n}$. The result also fails when the sets $\Omega$ are not assumed to be Stein; example: $X=$ blow-up of $\mathbb{C}^{n}$ at $0, T=[E]=$ current of integration on the exceptional divisor and $u(z)=\log |z|$ (see $\S 7.12$ for the definition of blow-ups).

Proof. By shrinking $\Omega$ slightly, we may assume that $\Omega$ has a smooth strongly pseudoconvex boundary. Let $\psi$ be a defining function of $\Omega$ as in (3.8). By subtracting a constant to $u$, we may assume $u \leqslant-\varepsilon$ on $\bar{\Omega}$. We fix $\delta$ so small that $\bar{\Omega} \backslash \Omega_{\delta}$ does not intersect $L(u) \cap$ Supp $T$ and we select a neighborhood $\omega$ of $\left(\bar{\Omega} \backslash \Omega_{\delta}\right) \cap$ Supp $T$ such that $\bar{\omega} \cap L(u)=\emptyset$. Then we define

$$
u_{s}(z)= \begin{cases}\max \{u(z), A \psi(z)\} & \text { on } \omega, \\ \max \{u(z), s\} & \text { on } \Omega_{\delta}=\{\psi<-\delta\} .\end{cases}
$$

By construction $u \geqslant-M$ on $\omega$ for some constant $M>0$. We fix $A \geqslant M / \delta$ and take $s \leqslant-M$, so

$$
\max \{u(z), A \psi(z)\}=\max \{u(z), s\}=u(z) \quad \text { on } \omega \cap \Omega_{\delta}
$$

and our definition of $u_{s}$ is coherent. Observe that $u_{s}$ is defined on $\omega \cup \Omega_{\delta}$, which is a neighborhood of $\bar{\Omega} \cap \operatorname{Supp} T$. Now, $u_{s}=A \psi$ on $\omega \cap\left(\Omega \backslash \Omega_{\varepsilon / A}\right)$, hence Stokes' theorem implies

$$
\begin{aligned}
\int_{\Omega} d d^{c} u_{s} \wedge T \wedge\left(d d^{c} \psi\right)^{p-1} & -\int_{\Omega} A d d^{c} \psi \wedge T \wedge\left(d d^{c} \psi\right)^{p-1} \\
& =\int_{\Omega} d d^{c}\left[\left(u_{s}-A \psi\right) T \wedge\left(d d^{c} \psi\right)^{p-1}\right]=0
\end{aligned}
$$

because the current [...] has a compact support contained in $\bar{\Omega}_{\varepsilon / A}$. Since $u_{s}$ and $\psi$ both vanish on $\partial \Omega$, an integration by parts gives

$$
\begin{aligned}
\int_{\Omega} u_{s} T \wedge\left(d d^{c} \psi\right)^{p} & =\int_{\Omega} \psi d d^{c} u_{s} \wedge T \wedge\left(d d^{c} \psi\right)^{p-1} \\
& \geqslant-\|\psi\|_{L^{\infty}(\Omega)} \int_{\Omega} T \wedge d d^{c} u_{s} \wedge\left(d d^{c} \psi\right)^{p-1} \\
& =-\|\psi\|_{L^{\infty}(\Omega)} A \int_{\Omega} T \wedge\left(d d^{c} \psi\right)^{p} .
\end{aligned}
$$

Finally, take $A=M / \delta$, let $s$ tend to $-\infty$ and use the inequality $u \geqslant-M$ on $\omega$. We obtain

$$
\begin{aligned}
\int_{\Omega} u T \wedge\left(d d^{c} \psi\right)^{p} & \geqslant-M \int_{\omega} T \wedge\left(d d^{c} \psi\right)^{p}+\lim _{s \rightarrow-\infty} \int_{\Omega_{\delta}} u_{s} T \wedge\left(d d^{c} \psi\right)^{p} \\
& \geqslant-\left(M+\|\psi\|_{L^{\infty}(\Omega)} M / \delta\right) \int_{\Omega} T \wedge\left(d d^{c} \psi\right)^{p}
\end{aligned}
$$

The last integral is finite. This concludes the proof.
(4.2) Remark. If $\Omega$ is smooth and strongly pseudoconvex, the above proof shows in fact that

$$
\left.\left\|\left.u T\right|_{\bar{\Omega}} \leqslant \frac{C}{\delta}\right\| u\left\|_{L^{\infty}\left(\left(\bar{\Omega} \backslash \Omega_{\delta}\right) \cap \operatorname{Supp} T\right)}\right\| T\right|_{\bar{\Omega}}
$$

when $L(u) \cap \operatorname{Supp} T \subset \Omega_{\delta}$. In fact, if $u$ is continuous and if $\omega$ is chosen sufficiently small, the constant $M$ can be taken arbitrarily close to $\|u\|_{L^{\infty}\left(\left(\bar{\Omega} \backslash \Omega_{\delta}\right) \cap \text { Supp } T\right)}$. Moreover, the maximum principle implies

$$
\left\|u_{+}\right\|_{L^{\infty}(\bar{\Omega} \cap \operatorname{Supp} T)}=\left\|u_{+}\right\|_{L^{\infty}(\partial \Omega \cap \operatorname{Supp} T)},
$$

so we can achieve that $u<-\varepsilon$ on a neighborhood of $\bar{\Omega} \cap \operatorname{Supp} T$ by subtracting $\|u\|_{L^{\infty}\left(\left(\bar{\Omega} \backslash \Omega_{\delta}\right) \cap \operatorname{Supp} T\right)}+2 \varepsilon$ [Proof of maximum principle: if $u\left(z_{0}\right)>0$ at $z_{0} \in \Omega \cap \operatorname{Supp} T$ and $u \leqslant 0$ near $\partial \Omega \cap \operatorname{Supp} T$, then

$$
\int_{\Omega} u_{+} T \wedge\left(d d^{c} \psi\right)^{p}=\int_{\Omega} \psi d d^{c} u_{+} \wedge T \wedge\left(d d^{c} \psi\right)^{p-1} \leqslant 0
$$

a contradiction].
(4.3) Corollary. Let $u_{1}, \ldots, u_{q}$ be plurisubharmonic functions on $X$ such that $X$ is covered by Stein open sets $\Omega$ with $\partial \Omega \cap L\left(u_{j}\right) \cap$ Supp $T=\emptyset$. We use again induction to define

$$
d d^{c} u_{1} \wedge d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T=d d^{c}\left(u_{1} d d^{c} u_{2} \ldots \wedge d d^{c} u_{q} \wedge T\right)
$$

Then, if $u_{1}^{k}, \ldots, u_{q}^{k}$ are decreasing sequences of plurisubharmonic functions converging pointwise to $u_{1}, \ldots, u_{q}, q \leqslant p$, properties ( $3.7 \mathrm{a}, \mathrm{b}$ ) hold.


III-2 Modified construction of $v_{j}^{k}$
Proof. Same proof as for Th. 3.7, with the following minor modification: the max procedure $v_{j}^{k}:=\max \left\{u_{j}^{k}, A \psi\right\}$ is applied only on a neighborhood $\omega$ of Supp $T \cap\left(\bar{\Omega} \backslash \Omega_{\delta}\right)$ with $\delta>0$ small, and $u_{j}^{k}$ is left unchanged near Supp $T \cap \bar{\Omega}_{\delta}$. Observe that the integration by part process requires the functions $u_{j}^{k}$ and $u_{j}^{k, \varepsilon}$ to be defined only near $\bar{\Omega} \cap \operatorname{Supp} T$.
(4.4) Proposition. Let $\Omega \subset \subset X$ be a Stein open subset. If $V$ is a plurisubharmonic function on $X$ and $u_{1}, \ldots, u_{q}, 1 \leqslant q \leqslant n-1$, are plurisubharmonic functions such that $\partial \Omega \cap L\left(u_{j}\right)=\emptyset$, then $V d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q}$ has locally finite mass in $\Omega$.

Proof. Same proof as for 3.11 , when $\delta>0$ is taken so small that $\Omega_{\delta} \supset L\left(u_{j}\right)$ for all $1 \leqslant j \leqslant q$.

Finally, we show that Monge-Ampère operators can also be defined in the case of plurisubharmonic functions with non compact pole sets, provided that the mutual intersections of the pole sets are of sufficiently small Hausdorff dimension with respect to the dimension $p$ of $T$.
(4.5) Theorem. Let $u_{1}, \ldots, u_{q}$ be plurisubharmonic functions on $X$. The currents $u_{1} d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T$ and $d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T$ are well defined and have locally finite mass in $X$ as soon as $q \leqslant p$ and

$$
\mathscr{H}_{2 p-2 m+1}\left(L\left(u_{j_{1}}\right) \cap \ldots \cap L\left(u_{j_{m}}\right) \cap \operatorname{Supp} T\right)=0
$$

for all choices of indices $j_{1}<\ldots<j_{m}$ in $\{1, \ldots, q\}$.

The proof is an easy induction on $q$, thanks to the following improved version of the Chern-Levine-Nirenberg inequalities.
(4.6) Proposition. Let $A_{1}, \ldots, A_{q} \subset X$ be closed sets such that

$$
\mathscr{H}_{2 p-2 m+1}\left(A_{j_{1}} \cap \ldots \cap A_{j_{m}} \cap \operatorname{Supp} T\right)=0
$$

for all choices of $j_{1}<\ldots<j_{m}$ in $\{1, \ldots, q\}$. Then for all compact sets $K, L$ of $X$ with $L \subset K^{\circ}$, there exist neighborhoods $V_{j}$ of $K \cap A_{j}$ and a constant $C=C\left(K, L, A_{j}\right)$ such that the conditions $u_{j} \leqslant 0$ on $K$ and $L\left(u_{j}\right) \subset A_{j}$ imply
a) $\left\|u_{1} d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T\right\|_{L} \leqslant C\left\|u_{1}\right\|_{L^{\infty}\left(K \backslash V_{1}\right)} \ldots\left\|u_{q}\right\|_{L^{\infty}\left(K \backslash V_{q}\right)}\|T\|_{K}$
b) $\left\|d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T\right\|_{L} \leqslant C\left\|u_{1}\right\|_{L^{\infty}\left(K \backslash V_{1}\right)} \ldots\left\|u_{q}\right\|_{L^{\infty}\left(K \backslash V_{q}\right)}\|T\|_{K}$.

Proof. We need only show that every point $x_{0} \in K^{\circ}$ has a neighborhood $L$ such that a), b) hold. Hence it is enough to work in a coordinate open set. We may thus assume that $X \subset \mathbb{C}^{n}$ is open, and after a regularization process $u_{j}=\lim u_{j} \star \rho_{\varepsilon}$ for $j=q, q-1, \ldots, 1$ in this order, that $u_{1}, \ldots, u_{q}$ are smooth. We proceed by induction on $q$ in two steps:

Step 1. $\left(\mathrm{b}_{q-1}\right) \Longrightarrow\left(\mathrm{b}_{q}\right)$,
Step 2. $\left(\mathrm{a}_{q-1}\right)$ and $\left(\mathrm{b}_{q}\right) \Longrightarrow\left(\mathrm{a}_{q}\right)$,
where $\left(\mathrm{b}_{0}\right)$ is the trivial statement $\|T\|_{L} \leqslant\|T\|_{K}$ and ( $\mathrm{a}_{0}$ ) is void. Observe that we have $\left(\mathrm{a}_{q}\right) \Longrightarrow\left(\mathrm{a}_{\ell}\right)$ and $\left(\mathrm{b}_{q}\right) \Longrightarrow\left(\mathrm{b}_{\ell}\right)$ for $\ell \leqslant q \leqslant p$ by taking $u_{\ell+1}(z)=\ldots=u_{q}(z)=|z|^{2}$. We need the following elementary fact.
(4.7) Lemma. Let $F \subset \mathbb{C}^{n}$ be a closed set such that $\mathscr{H}_{2 s+1}(F)=0$ for some integer $0 \leqslant s<n$. Then for almost all choices of unitary coordinates $\left(z_{1}, \ldots, z_{n}\right)=\left(z^{\prime}, z^{\prime \prime}\right)$ with $z^{\prime}=\left(z_{1}, \ldots, z_{s}\right), z^{\prime \prime}=\left(z_{s+1}, \ldots, z_{n}\right)$ and almost all radii of balls $B^{\prime \prime}=B\left(0, r^{\prime \prime}\right) \subset \mathbb{C}^{n-s}$, the set $\{0\} \times \partial B^{\prime \prime}$ does not intersect $F$.

Proof. The unitary group $U(n)$ has real dimension $n^{2}$. There is a proper submersion

$$
\Phi: U(n) \times\left(\mathbb{C}^{n-s} \backslash\{0\}\right) \longrightarrow \mathbb{C}^{n} \backslash\{0\}, \quad\left(g, z^{\prime \prime}\right) \longmapsto g\left(0, z^{\prime \prime}\right)
$$

whose fibers have real dimension $N=n^{2}-2 s$. It follows that the inverse image $\Phi^{-1}(F)$ has zero Hausdorff measure $\mathscr{H}_{N+2 s+1}=\mathscr{H}_{n^{2}+1}$. The set of pairs $\left(g, r^{\prime \prime}\right) \in U(n) \times \mathbb{R}_{+}^{\star}$ such that $g\left(\{0\} \times \partial B^{\prime \prime}\right)$ intersects $F$ is precisely the image of $\Phi^{-1}(F)$ in $U(n) \times \mathbb{R}_{+}^{\star}$ by the Lipschitz map $\left(g, z^{\prime \prime}\right) \mapsto\left(g,\left|z^{\prime \prime}\right|\right)$. Hence this set has zero $\mathscr{H}_{n^{2}+1}$-measure.

Proof of step 1.. Take $x_{0}=0 \in K^{\circ}$. Suppose first $0 \in A_{1} \cap \ldots \cap A_{q}$ and set $F=$ $A_{1} \cap \ldots \cap A_{q} \cap$ Supp $T$. Since $\mathscr{H}_{2 p-2 q+1}(F)=0$, Lemma 4.7 implies that there are coordinates $z^{\prime}=\left(z_{1}, \ldots, z_{s}\right), z^{\prime \prime}=\left(z_{s+1}, \ldots, z_{n}\right)$ with $s=p-q$ and a ball $\bar{B}^{\prime \prime}$ such that $F \cap\left(\{0\} \times \partial B^{\prime \prime}\right)=\emptyset$ and $\{0\} \times \bar{B}^{\prime \prime} \subset K^{\circ}$. By compactness of $K$, we can find neighborhoods $W_{j}$ of $K \cap A_{j}$ and a ball $B^{\prime}=B\left(0, r^{\prime}\right) \subset \mathbb{C}^{s}$ such that $\bar{B}^{\prime} \times \bar{B}^{\prime \prime} \subset K^{\circ}$ and

$$
\begin{equation*}
\bar{W}_{1} \cap \ldots \cap \bar{W}_{q} \cap \operatorname{Supp} T \cap\left(\bar{B}^{\prime} \times\left(\bar{B}^{\prime \prime} \backslash(1-\delta) B^{\prime \prime}\right)\right)=\emptyset \tag{4.8}
\end{equation*}
$$

for $\delta>0$ small. If $0 \notin A_{j}$ for some $j$, we choose instead $W_{j}$ to be a small neighborhood of 0 such that $\bar{W}_{j} \subset\left(\bar{B}^{\prime} \times(1-\delta) B^{\prime \prime}\right) \backslash A_{j} ;$ property (4.8) is then automatically satisfied.

Let $\chi_{j} \geqslant 0$ be a function with compact support in $W_{j}$, equal to 1 near $K \cap A_{j}$ if $A_{j} \ni 0$ (resp. equal to 1 near 0 if $A_{j} \not \supset 0$ ) and let $\chi\left(z^{\prime}\right) \geqslant 0$ be a function equal to 1 on $1 / 2 B^{\prime}$ with compact support in $B^{\prime}$. Then

$$
\int_{B^{\prime} \times B^{\prime \prime}} d d^{c}\left(\chi_{1} u_{1}\right) \wedge \ldots \wedge d d^{c}\left(\chi_{q} u_{q}\right) \wedge T \wedge \chi\left(z^{\prime}\right)\left(d d^{c}\left|z^{\prime}\right|^{2}\right)^{s}=0
$$

because the integrand is $d d^{c}$ exact and has compact support in $B^{\prime} \times B^{\prime \prime}$ thanks to (4.8). If we expand all factors $d d^{c}\left(\chi_{j} u_{j}\right)$, we find a term

$$
\chi_{1} \ldots \chi_{q} \chi\left(z^{\prime}\right) d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \geqslant 0
$$

which coincides with $d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T$ on a small neighborhood of 0 where $\chi_{j}=\chi=1$. The other terms involve

$$
d \chi_{j} \wedge d^{c} u_{j}+d u_{j} \wedge d^{c} \chi_{j}+u_{j} d d^{c} \chi_{j}
$$

for at least one index $j$. However $d \chi_{j}$ and $d d^{c} \chi^{\prime}$, vanish on some neighborhood $V_{j}^{\prime}$ of $K \cap A_{j}$ and therefore $u_{j}$ is bounded on $\bar{B}^{\prime} \times \bar{B}^{\prime \prime \prime} \backslash V_{j}^{\prime}$. We then apply the induction hypothesis $\left(\mathrm{b}_{q-1}\right)$ to the current

$$
\Theta=d d^{c} u_{1} \wedge \ldots \wedge \widehat{d d^{c} u_{j}} \wedge \ldots \wedge d d^{c} u_{q} \wedge T
$$

and the usual Chern-Levine-Nirenberg inequality to the product of $\Theta$ with the mixed term $d \chi_{j} \wedge d^{c} u_{j}+d u_{j} \wedge d^{c} \chi_{j}$. Remark 3.5 can be applied because $\chi_{j}$ is smooth and is therefore a difference $\chi_{j}^{(1)}-\chi_{j}^{(2)}$ of locally bounded plurisubharmonic functions in $\mathbb{C}^{n}$. Let $K^{\prime}$ be a compact neighborhood of $\bar{B}^{\prime} \times \bar{B}^{\prime \prime}$ with $K_{\bar{B}^{\prime \prime}}^{\prime \prime} \subset K^{\circ}$, and let $V_{j}$ be a neighborhood of $K \cap A_{j}$ with $\bar{V}_{j} \subset V_{j}^{\prime}$. Then with $L^{\prime}:=\left(\bar{B}^{\prime} \times \bar{B}^{\prime \prime}\right) \backslash V_{j}^{\prime} \subset\left(K^{\prime} \backslash V_{j}\right)^{\circ}$ we obtain

$$
\begin{aligned}
& \|\left(d \chi_{j} \wedge d^{c} u_{j}+d u_{j} \wedge d^{c} \chi_{j}\right) \wedge \Theta\left\|_{\bar{B}^{\prime} \times \bar{B}^{\prime \prime}}=\right\|\left(d \chi_{j} \wedge d^{c} u_{j}+d u_{j} \wedge d^{c} \chi_{j}\right) \wedge \Theta \|_{L^{\prime}} \\
& \leqslant C_{1}\left\|u_{j}\right\|_{L^{\infty}\left(K^{\prime} \backslash V_{j}\right)}\|\Theta\|_{K^{\prime} \backslash V_{j}} \\
&\|\Theta\|_{K^{\prime} \backslash V_{j}} \leqslant\|\Theta\|_{K^{\prime}} \leqslant C_{2}\left\|u_{1}\right\|_{L^{\infty}\left(K \backslash V_{1}\right)} \ldots\left\|u_{j}\right\| \ldots\left\|u_{q}\right\|_{L^{\infty}\left(K \backslash V_{q}\right)}\|T\|_{K} .
\end{aligned}
$$

Now, we may slightly move the unitary basis in $\mathbb{C}^{n}$ and get coordinate systems $z^{m}=$ $\left(z_{1}^{m}, \ldots, z_{n}^{m}\right)$ with the same properties as above, such that the forms

$$
\left(d d^{c}\left|z^{m \prime}\right|^{2}\right)^{s}=\frac{s!}{\pi^{s}} \mathrm{i} d z_{1}^{m} \wedge d \bar{z}_{1}^{m} \wedge \ldots \wedge \mathrm{i} d z_{s}^{m} \wedge d \bar{z}_{s}^{m}, \quad 1 \leqslant m \leqslant N
$$

define a basis of $\bigwedge^{s, s}\left(\mathbb{C}^{n}\right)^{\star}$. It follows that all measures

$$
d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \wedge \mathrm{i} d z_{1}^{m} \wedge d \bar{z}_{1}^{m} \wedge \ldots \wedge \mathrm{i} d z_{s}^{m} \wedge d \bar{z}_{s}^{m}
$$

satisfy estimate $\left(\mathrm{b}_{q}\right)$ on a small neighborhood $L$ of 0 .
Proof of Step 2.. We argue in a similar way with the integrals

$$
\begin{aligned}
\int_{B^{\prime} \times B^{\prime \prime}} & \chi_{1} u_{1} d d^{c}\left(\chi_{2} u_{2}\right) \wedge \ldots d d^{c}\left(\chi_{q} u_{q}\right) \wedge T \wedge \chi\left(z^{\prime}\right)\left(d d^{c}\left|z^{\prime}\right|^{2}\right)^{s} \wedge d d^{c}\left|z_{s+1}\right|^{2} \\
& =\int_{B^{\prime} \times B^{\prime \prime}}\left|z_{s+1}\right|^{2} d d^{c}\left(\chi_{1} u_{1}\right) \wedge \ldots d d^{c}\left(\chi_{q} u_{q}\right) \wedge T \wedge \chi\left(z^{\prime}\right)\left(d d^{c}\left|z^{\prime}\right|^{2}\right)^{s}
\end{aligned}
$$

We already know by $\left(\mathrm{b}_{q}\right)$ and Remark 3.5 that all terms in the right hand integral admit the desired bound. For $q=1$, this shows that $\left(\mathrm{b}_{1}\right) \Longrightarrow\left(\mathrm{a}_{1}\right)$. Except for $\chi_{1} \ldots \chi_{q} \chi\left(z^{\prime}\right) u_{1} d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T$, all terms in the left hand integral involve derivatives of $\chi_{j}$. By construction, the support of these derivatives is disjoint from $A_{j}$, thus we only have to obtain a bound for

$$
\int_{L} u_{1} d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \wedge \alpha
$$

when $L=\bar{B}\left(x_{0}, r\right)$ is disjoint from $A_{j}$ for some $j \geqslant 2$, say $L \cap A_{2}=\emptyset$, and $\alpha$ is a constant positive form of type $(p-q, p-q)$. Then $\bar{B}\left(x_{0}, r+\varepsilon\right) \subset K^{\circ} \backslash \bar{V}_{2}$ for some $\varepsilon>0$ and some neighborhood $V_{2}$ of $K \cap A_{2}$. By the max construction used e.g. in Prop. 4.1, we can replace $u_{2}$ by a plurisubharmonic function $\widetilde{u}_{2}$ equal to $u_{2}$ in $L$ and to $A\left(\left|z-x_{0}\right|^{2}-r^{2}\right)-M$ in $\bar{B}\left(x_{0}, r+\varepsilon\right) \backslash B\left(x_{0}, r+\varepsilon / 2\right)$, with $M=\left\|u_{2}\right\|_{L^{\infty}\left(K \backslash V_{2}\right)}$ and $A=M / \varepsilon r$. Let $\chi \geqslant 0$ be a smooth function equal to 1 on $B\left(x_{0}, r+\varepsilon / 2\right)$ with support in $B\left(x_{0}, r\right)$. Then

$$
\begin{array}{rl}
\int_{B\left(x_{0}, r+\varepsilon\right)} u_{1} & d d^{c}\left(\chi \widetilde{u}_{2}\right) \wedge d d^{c} u_{3} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \wedge \alpha \\
& =\int_{B\left(x_{0}, r+\varepsilon\right)} \chi \widetilde{u}_{2} d d^{c} u_{1} \wedge d d^{c} u_{3} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \wedge \alpha \\
& \leqslant O(1)\left\|u_{1}\right\|_{L^{\infty}\left(K \backslash V_{1}\right)} \ldots\left\|u_{q}\right\|_{L^{\infty}\left(K \backslash V_{q}\right)}\|T\|_{K}
\end{array}
$$

where the last estimate is obtained by the induction hypothesis $\left(\mathrm{b}_{q-1}\right)$ applied to $d d^{c} u_{1} \wedge$ $d d^{c} u_{3} \wedge \ldots \wedge d d^{c} u_{q} \wedge T$. By construction

$$
d d^{c}\left(\chi \widetilde{u}_{2}\right)=\chi d d^{c} \widetilde{u}_{2}+(\text { smooth terms involving } d \chi)
$$

coincides with $d d^{c} u_{2}$ in $L$, and ( $\mathrm{a}_{q-1}$ ) implies the required estimate for the other terms in the left hand integral.
(4.9) Proposition. With the assumptions of Th. 4.5, the analogue of the monotone convergence Theorem 3.7 (a,b) holds.

Proof. By the arguments already used in the proof of Th. 3.7 (e.g. by Lemma 3.9), it is enough to show that

$$
\begin{aligned}
\int_{B^{\prime} \times B^{\prime \prime}} & \chi_{1} \ldots \chi_{q} u_{1} \wedge d d^{c} u_{2} \wedge \ldots \wedge d d^{c} u_{q} \wedge T \wedge \alpha \\
& \leqslant \liminf _{k \rightarrow+\infty} \int_{B^{\prime} \times B^{\prime \prime}} \chi_{1} \ldots \chi_{q} u_{1}^{k} d d^{c} u_{2}^{k} \wedge \ldots \wedge d d^{c} u_{q}^{k} \wedge T \wedge \alpha
\end{aligned}
$$

where $\alpha=\chi\left(z^{\prime}\right)\left(d d^{c}\left|z^{\prime}\right|^{2}\right)^{s}$ is closed. Here the functions $\chi_{j}, \chi$ are chosen as in the proof of Step 1 in 4.7 , especially their product has compact support in $B^{\prime} \times B^{\prime \prime}$ and $\chi_{j}=\chi=1$ in a neighborhood of the given point $x_{0}$. We argue by induction on $q$ and also on the number $m$ of functions $\left(u_{j}\right)_{j \geqslant 1}$ which are unbounded near $x_{0}$. If $u_{j}$ is bounded near $x_{0}$, we take $W_{j}^{\prime \prime} \subset \subset W_{j}^{\prime} \subset \subset W_{j}$ to be small balls of center $x_{0}$ on which $u_{j}$ is bounded and we modify the sequence $u_{j}^{k}$ on the corona $W_{j} \backslash W_{j}^{\prime \prime}$ so as to make it constant and equal to a smooth function $A\left|z-x_{0}\right|^{2}+B$ on the smaller corona $W_{j} \backslash W_{j}^{\prime}$. In that case, we
take $\chi_{j}$ equal to 1 near $\bar{W}_{j}^{\prime}$ and Supp $\chi_{j} \subset W_{j}$. For every $\ell=1, \ldots, q$, we are going to check that

$$
\begin{aligned}
& \quad \liminf _{k \rightarrow+\infty} \int_{B^{\prime} \times B^{\prime \prime}} \chi_{1} u_{1}^{k} d d^{c}\left(\chi_{2} u_{2}^{k}\right) \wedge \ldots \\
& \quad d d^{c}\left(\chi_{\ell-1} u_{\ell-1}^{k}\right) \wedge d d^{c}\left(\chi_{\ell} u_{\ell}\right) \wedge d d^{c}\left(\chi_{\ell+1} u_{\ell+1}\right) \ldots d d^{c}\left(\chi_{q} u_{q}\right) \wedge T \wedge \alpha \\
& \leqslant \liminf _{k \rightarrow+\infty} \int_{B^{\prime} \times B^{\prime \prime}} \chi_{1} u_{1}^{k} d d^{c}\left(\chi_{2} u_{2}^{k}\right) \wedge \ldots \\
& \quad d d^{c}\left(\chi_{\ell-1} u_{\ell-1}^{k}\right) \wedge d d^{c}\left(\chi_{\ell} u_{\ell}^{k}\right) \wedge d d^{c}\left(\chi_{\ell+1} u_{\ell+1}\right) \ldots d d^{c}\left(\chi_{q} u_{q}\right) \wedge T \wedge \alpha .
\end{aligned}
$$

In order to do this, we integrate by parts $\chi_{1} u_{1}^{k} d d^{c}\left(\chi_{\ell} u_{\ell}\right)$ into $\chi_{\ell} u_{\ell} d d^{c}\left(\chi_{1} u_{1}^{k}\right)$ for $\ell \geqslant 2$, and we use the inequality $u_{\ell} \leqslant u_{\ell}^{k}$. Of course, the derivatives $d \chi_{j}, d^{c} \chi_{j}, d d^{c} \chi_{j}$ produce terms which are no longer positive and we have to take care of these. However, Supp $d \chi_{j}$ is disjoint from the unbounded locus of $u_{j}$ when $u_{j}$ is unbounded, and contained in $W_{j} \backslash \bar{W}_{j}^{\prime}$ when $u_{j}$ is bounded. The number $m$ of unbounded functions is therefore replaced by $m-1$ in the first case, whereas in the second case $u_{j}^{k}=u_{j}$ is constant and smooth on Supp $d \chi_{j}$, so $q$ can be replaced by $q-1$. By induction on $q+m$ (and thanks to the polarization technique 3.5), the limit of the terms involving derivatives of $\chi_{j}$ is equal on both sides to the corresponding terms obtained by suppressing all indices $k$. Hence these terms do not give any contribution in the inequalities.

We finally quote the following simple consequences of Th. 4.5 when $T$ is arbitrary and $q=1$, resp. when $T=1$ has bidegree $(0,0)$ and $q$ is arbitrary.
(4.10) Corollary. Let $T$ be a closed positive current of bidimension $(p, p)$ and let $u$ be a plurisubharmonic function on $X$ such that $L(u) \cap$ Supp $T$ is contained in an analytic set of dimension at most $p-1$. Then $u T$ and $d d^{c} u \wedge T$ are well defined and have locally finite mass in $X$.
(4.11) Corollary. Let $u_{1}, \ldots, u_{q}$ be plurisubharmonic functions on $X$ such that $L\left(u_{j}\right)$ is contained in an analytic set $A_{j} \subset X$ for every $j$. Then $d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q}$ is well defined as soon as $A_{j_{1}} \cap \ldots \cap A_{j_{m}}$ has codimension at least $m$ for all choices of indices $j_{1}<\ldots<j_{m}$ in $\{1, \ldots, q\}$.

In the particular case when $u_{j}=\log \left|f_{j}\right|$ for some non zero holomorphic function $f_{j}$ on $X$, we see that the intersection product of the associated zero divisors $\left[Z_{j}\right]=d d^{c} u_{j}$ is well defined as soon as the supports $\left|Z_{j}\right|$ satisfy $\operatorname{codim}\left|Z_{j_{1}}\right| \cap \ldots \cap\left|Z_{j_{m}}\right|=m$ for every $m$. Similarly, when $T=[A]$ is an analytic $p$-cycle, Cor. 4.10 shows that $[Z] \wedge[A]$ is well defined for every divisor $Z$ such that $\operatorname{dim}|Z| \cap|A|=p-1$. These observations easily imply the following
(4.12) Proposition. Suppose that the divisors $Z_{j}$ satisfy the above codimension condition and let $\left(C_{k}\right)_{k \geqslant 1}$ be the irreducible components of the point set intersection $\left|Z_{1}\right| \cap$ $\ldots \cap\left|Z_{q}\right|$. Then there exist integers $m_{k}>0$ such that

$$
\left[Z_{1}\right] \wedge \ldots \wedge\left[Z_{q}\right]=\sum m_{k}\left[C_{k}\right] .
$$

The integer $m_{k}$ is called the multiplicity of intersection of $Z_{1}, \ldots, Z_{q}$ along the component $C_{k}$.

Proof. The wedge product has bidegree $(q, q)$ and support in $C=\bigcup C_{k}$ where codim $C=$ $q$, so it must be a sum as above with $m_{k} \in \mathbb{R}_{+}$. We check by induction on $q$ that $m_{k}$ is a positive integer. If we denote by $A$ some irreducible component of $\left|Z_{1}\right| \cap \ldots \cap\left|Z_{q-1}\right|$, we need only check that $[A] \wedge\left[Z_{q}\right]$ is an integral analytic cycle of codimension $q$ with positive coefficients on each component $C_{k}$ of the intersection. However $[A] \wedge\left[Z_{q}\right]=$ $d d^{c}\left(\log \left|f_{q}\right|[A]\right)$. First suppose that no component of $A \cap f_{q}^{-1}(0)$ is contained in the singular part $A_{\text {sing }}$. Then the Lelong-Poincaré equation applied on $A_{\text {reg }}$ shows that $d d^{c}\left(\log \left|f_{q}\right|[A]\right)=\sum m_{k}\left[C_{k}\right]$ on $X \backslash A_{\text {sing }}$, where $m_{k}$ is the vanishing order of $f_{q}$ along $C_{k}$ in $A_{\text {reg. }}$. Since $C \cap A_{\text {sing }}$ has codimension $q+1$ at least, the equality must hold on $X$. In general, we replace $f_{q}$ by $f_{q}-\varepsilon$ so that the divisor of $f_{q}-\varepsilon$ has no component contained in $A_{\text {sing }}$. Then $d d^{c}\left(\log \left|f_{q}-\varepsilon\right|[A]\right)$ is an integral codimension $q$ cycle with positive multiplicities on each component of $A \cap f_{q}^{-1}(\varepsilon)$ and we conclude by letting $\varepsilon$ tend to zero.

## § 5. Generalized Lelong Numbers

The concepts we are going to study mostly concern the behaviour of currents or plurisubharmonic functions in a neighborhood of a point at which we have for instance a logarithmic pole. Since the interesting applications are local, we assume from now on (unless otherwise stated) that $X$ is a Stein manifold, i.e. that $X$ has a strictly plurisubharmonic exhaustion function. Let $\varphi: X \longrightarrow[-\infty,+\infty[$ be a continuous plurisubharmonic function (in general $\varphi$ may have $-\infty$ poles, our continuity assumption means that $e^{\varphi}$ is continuous). The sets

$$
\begin{align*}
S(r) & =\{x \in X ; \varphi(x)=r\},  \tag{5.1}\\
B(r) & =\{x \in X ; \varphi(x)<r\}, \\
\bar{B}(r) & =\{x \in X ; \varphi(x) \leqslant r\}
\end{align*}
$$

will be called pseudo-spheres and pseudo-balls associated with $\varphi$. Note that $\bar{B}(r)$ is not necessarily equal to the closure of $B(r)$, but this is often true in concrete situations. The most simple example we have in mind is the case of the function $\varphi(z)=\log |z-a|$ on an open subset $X \subset \mathbb{C}^{n}$; in this case $B(r)$ is the euclidean ball of center $a$ and radius $e^{r}$; moreover, the forms

$$
\begin{equation*}
\frac{1}{2} d d^{c} e^{2 \varphi}=\frac{\mathrm{i}}{2 \pi} d^{\prime} d^{\prime \prime}|z|^{2}, \quad d d^{c} \varphi=\frac{\mathrm{i}}{\pi} d^{\prime} d^{\prime \prime} \log |z-a| \tag{5.2}
\end{equation*}
$$

can be interpreted respectively as the flat hermitian metric on $\mathbb{C}^{n}$ and as the pull-back over $\mathbb{C}^{n}$ of the Fubini-Study metric of $\mathbb{P}^{n-1}$, translated by $a$.
(5.3) Definition. We say that $\varphi$ is semi-exhaustive if there exists a real number $R$ such that $B(R) \subset \subset X$. Similarly, $\varphi$ is said to be semi-exhaustive on a closed subset $A \subset X$ if there exists $R$ such that $A \cap B(R) \subset \subset X$.

We are interested especially in the set of poles $S(-\infty)=\{\varphi=-\infty\}$ and in the behaviour of $\varphi$ near $S(-\infty)$. Let $T$ be a closed positive current of bidimension $(p, p)$ on $X$. Assume that $\varphi$ is semi-exhaustive on Supp $T$ and that $B(R) \cap \operatorname{Supp} T \subset \subset X$. Then
$P=S(-\infty) \cap$ SuppT is compact and the results of $\S 2$ show that the measure $T \wedge\left(d d^{c} \varphi\right)^{p}$ is well defined. Following [Demailly 1982b, 1987a], we introduce:
(5.4) Definition. If $\varphi$ is semi-exhaustive on Supp $T$ and if $R$ is such that $B(R) \cap$ Supp $T \subset \subset$, we set for all $r \in]-\infty, R[$

$$
\begin{aligned}
\nu(T, \varphi, r) & =\int_{B(r)} T \wedge\left(d d^{c} \varphi\right)^{p}, \\
\nu(T, \varphi) & =\int_{S(-\infty)} T \wedge\left(d d^{c} \varphi\right)^{p}=\lim _{r \rightarrow-\infty} \nu(T, \varphi, r) .
\end{aligned}
$$

The number $\nu(T, \varphi)$ will be called the (generalized) Lelong number of $T$ with respect to the weight $\varphi$.

If we had not required $T \wedge\left(d d^{c} \varphi\right)^{p}$ to be defined pointwise on $\varphi^{-1}(-\infty)$, the assumption that $X$ is Stein could have been dropped: in fact, the integral over $B(r)$ always makes sense if we define

$$
\nu(T, \varphi, r)=\int_{B(r)} T \wedge\left(d d^{c} \max \{\varphi, s\}\right)^{p} \quad \text { with } s<r
$$

Stokes' formula shows that the right hand integral is actually independent of $s$. The example given after (4.1) shows however that $T \wedge\left(d d^{c} \varphi\right)^{p}$ need not exist on $\varphi^{-1}(-\infty)$ if $\varphi^{-1}(-\infty)$ contains an exceptional compact analytic subset. We leave the reader consider by himself this more general situation and extend our statements by the $\max \{\varphi, s\}$ technique. Observe that $r \longmapsto \nu(T, \varphi, r)$ is always an increasing function of $r$. Before giving examples, we need a formula.
(5.5) Formula. For any convex increasing function $\chi: \mathbb{R} \longrightarrow \mathbb{R}$ we have

$$
\int_{B(r)} T \wedge\left(d d^{c} \chi \circ \varphi\right)^{p}=\chi^{\prime}(r-0)^{p} \nu(T, \varphi, r)
$$

where $\chi^{\prime}(r-0)$ denotes the left derivative of $\chi$ at $r$.
Proof. Let $\chi_{\varepsilon}$ be the convex function equal to $\chi$ on $[r-\varepsilon,+\infty[$ and to a linear function of slope $\chi^{\prime}(r-\varepsilon-0)$ on $\left.]-\infty, r-\varepsilon\right]$. We get $d d^{c}\left(\chi_{\varepsilon} \circ \varphi\right)=\chi^{\prime}(r-\varepsilon-0) d d^{c} \varphi$ on $B(r-\varepsilon)$ and Stokes' theorem implies

$$
\begin{aligned}
\int_{B(r)} T \wedge\left(d d^{c} \chi \circ \varphi\right)^{p} & =\int_{B(r)} T \wedge\left(d d^{c} \chi_{\varepsilon} \circ \varphi\right)^{p} \\
& \geqslant \int_{B(r-\varepsilon)} T \wedge\left(d d^{c} \chi_{\varepsilon} \circ \varphi\right)^{p} \\
& =\chi^{\prime}(r-\varepsilon-0)^{p} \nu(T, \varphi, r-\varepsilon)
\end{aligned}
$$

Similarly, taking $\widetilde{\chi}_{\varepsilon}$ equal to $\chi$ on $\left.]-\infty, r-\varepsilon\right]$ and linear on $[r-\varepsilon, r]$, we obtain

$$
\int_{B(r-\varepsilon)} T \wedge\left(d d^{c} \chi \circ \varphi\right)^{p} \leqslant \int_{B(r)} T \wedge\left(d d^{c} \widetilde{\chi}_{\varepsilon} \circ \varphi\right)^{p}=\chi^{\prime}(r-\varepsilon-0)^{p} \nu(T, \varphi, r)
$$

The expected formula follows when $\varepsilon$ tends to 0 .
We get in particular $\int_{B(r)} T \wedge\left(d d^{c} e^{2 \varphi}\right)^{p}=\left(2 e^{2 r}\right)^{p} \nu(T, \varphi, r)$, whence the formula

$$
\begin{equation*}
\nu(T, \varphi, r)=e^{-2 p r} \int_{B(r)} T \wedge\left(\frac{1}{2} d d^{c} e^{2 \varphi}\right)^{p} . \tag{5.6}
\end{equation*}
$$

Now, assume that $X$ is an open subset of $\mathbb{C}^{n}$ and that $\varphi(z)=\log |z-a|$ for some $a \in X$. Formula (5.6) gives

$$
\nu(T, \varphi, \log r)=r^{-2 p} \int_{|z-a|<r} T \wedge\left(\frac{\mathrm{i}}{2 \pi} d^{\prime} d^{\prime \prime}|z|^{2}\right)^{p} .
$$

The positive measure $\sigma_{T}=\frac{1}{p!} T \wedge\left(\frac{\mathrm{i}}{2} d^{\prime} d^{\prime \prime}|z|^{2}\right)^{p}=2^{-p} \sum T_{I, I} \cdot \mathrm{i}^{n} d z_{1} \wedge \ldots \wedge d \bar{z}_{n}$ is called the trace measure of $T$. We get

$$
\begin{equation*}
\nu(T, \varphi, \log r)=\frac{\sigma_{T}(B(a, r))}{\pi^{p} r^{2 p} / p!} \tag{5.7}
\end{equation*}
$$

and $\nu(T, \varphi)$ is the limit of this ratio as $r \rightarrow 0$. This limit is called the (ordinary) Lelong number of $T$ at point $a$ and is denoted $\nu(T, a)$. This was precisely the original definition of Lelong, see [Lelong 1968]. Let us mention a simple but important consequence.
(5.8) Consequence. The ratio $\sigma_{T}(B(a, r)) / r^{2 p}$ is an increasing function of $r$. Moreover, for every compact subset $K \subset X$ and every $r_{0}<d(K, \partial X)$ we have

$$
\sigma_{T}(B(a, r)) \leqslant C r^{2 p} \quad \text { for } a \in K \text { and } r \leqslant r_{0}
$$

where $C=\sigma_{T}\left(K+\bar{B}\left(0, r_{0}\right)\right) / r_{0}^{2 p}$.

All these results are particularly interesting when $T=[A]$ is the current of integration over an analytic subset $A \subset X$ of pure dimension $p$. Then $\sigma_{T}(B(a, r))$ is the euclidean area of $A \cap B(a, r)$, while $\pi^{p} r^{2 p} / p!$ is the area of a ball of radius $r$ in a $p$-dimensional subspace of $\mathbb{C}^{n}$. Thus $\nu(T, \varphi, \log r)$ is the ratio of these areas and the Lelong number $\nu(T, a)$ is the limit ratio.
(5.9) Remark. It is immediate to check that

$$
\nu([A], x)= \begin{cases}0 & \text { for } x \notin A, \\ 1 & \text { when } x \in A \text { is a regular point. }\end{cases}
$$

We will see later that $\nu([A], x)$ is always an integer (Thie's theorem 8.7).
(5.10) Remark. When $X=\mathbb{C}^{n}, \varphi(z)=\log |z-a|$ and $A=X$ (i.e. $T=1$ ), we obtain in particular $\int_{B(a, r)}\left(d d^{c} \log |z-a|\right)^{n}=1$ for all $r$. This implies

$$
\left(d d^{c} \log |z-a|\right)^{n}=\delta_{a} .
$$

This fundamental formula can be viewed as a higher dimensional analogue of the usual formula $\Delta \log |z-a|=2 \pi \delta_{a}$ in $\mathbb{C}$.

We next prove a result which shows in particular that the Lelong numbers of a closed positive current are zero except on a very small set.
(5.11) Proposition. If $T$ is a closed positive current of bidimension ( $p, p$ ), then for each $c>0$ the set $E_{c}=\{x \in X ; \nu(T, x) \geqslant c\}$ is a closed set of locally finite $\mathscr{H}_{2 p}$ Hausdorff measure in $X$.

Proof. By (5.7), we infer $\nu(T, a)=\lim _{r \rightarrow 0} \sigma_{T}(\bar{B}(a, r)) p!/ \pi^{p} r^{2 p}$. The function $a \mapsto$ $\sigma_{T}(\bar{B}(a, r))$ is clearly upper semicontinuous. Hence the decreasing limit $\nu(T, a)$ as $r$ decreases to 0 is also upper semicontinuous in $a$. This implies that $E_{c}$ is closed. Now, let $K$ be a compact subset in $X$ and let $\left\{a_{j}\right\}_{1 \leqslant j \leqslant N}, N=N(\varepsilon)$, be a maximal collection of points in $E_{c} \cap K$ such that $\left|a_{j}-a_{k}\right| \geqslant 2 \varepsilon$ for $j \neq k$. The balls $B\left(a_{j}, 2 \varepsilon\right)$ cover $E_{c} \cap K$, whereas the balls $B\left(a_{j}, \varepsilon\right)$ are disjoint. If $K_{c, \varepsilon}$ is the set of points which are at distance $\leqslant \varepsilon$ of $E_{c} \cap K$, we get

$$
\sigma_{T}\left(K_{c, \varepsilon}\right) \geqslant \sum \sigma_{T}\left(B\left(a_{j}, \varepsilon\right)\right) \geqslant N(\varepsilon) c \pi^{p} \varepsilon^{2 p} / p!
$$

since $\nu\left(T, a_{j}\right) \geqslant c$. By the definition of Hausdorff measure, we infer

$$
\begin{aligned}
\mathscr{H}_{2 p}\left(E_{c} \cap K\right) & \leqslant \liminf _{\varepsilon \rightarrow 0} \sum\left(\operatorname{diam} B\left(a_{j}, 2 \varepsilon\right)\right)^{2 p} \\
& \leqslant \liminf _{\varepsilon \rightarrow 0} N(\varepsilon)(4 \varepsilon)^{2 p} \leqslant \frac{p!4^{2 p}}{c \pi^{p}} \sigma_{T}\left(E_{c} \cap K\right) .
\end{aligned}
$$

Finally, we conclude this section by proving two simple semi-continuity results for Lelong numbers.
(5.12) Proposition. Let $T_{k}$ be a sequence of closed positive currents of bidimension $(p, p)$ converging weakly to a limit $T$. Suppose that there is a closed set $A$ such that Supp $T_{k} \subset A$ for all $k$ and such that $\varphi$ is semi-exhaustive on $A$ with $A \cap B(R) \subset \subset$. Then for all $r<R$ we have

$$
\begin{aligned}
\int_{B(r)} T \wedge\left(d d^{c} \varphi\right)^{p} & \leqslant \liminf _{k \rightarrow+\infty} \int_{B(r)} T_{k} \wedge\left(d d^{c} \varphi\right)^{p} \\
& \leqslant \limsup _{k \rightarrow+\infty} \int_{\bar{B}(r)} T_{k} \wedge\left(d d^{c} \varphi\right)^{p} \leqslant \int_{\bar{B}(r)} T \wedge\left(d d^{c} \varphi\right)^{p}
\end{aligned}
$$

When $r$ tends to $-\infty$, we find in particular

$$
\limsup _{k \rightarrow+\infty} \nu\left(T_{k}, \varphi\right) \leqslant \nu(T, \varphi)
$$

Proof. Let us prove for instance the third inequality. Let $\varphi_{\ell}$ be a sequence of smooth plurisubharmonic approximations of $\varphi$ with $\varphi \leqslant \varphi_{\ell}<\varphi+1 / \ell$ on $\{r-\varepsilon \leqslant \varphi \leqslant r+\varepsilon\}$. We set

$$
\psi_{\ell}= \begin{cases}\varphi & \text { on } \bar{B}(r), \\ \max \left\{\varphi,(1+\varepsilon)\left(\varphi_{\ell}-1 / \ell\right)-r \varepsilon\right\} & \text { on } X \backslash B(r) .\end{cases}
$$

This definition is coherent since $\psi_{\ell}=\varphi$ near $S(r)$, and we have

$$
\psi_{\ell}=(1+\varepsilon)\left(\varphi_{\ell}-1 / \ell\right)-r \varepsilon \quad \text { near } \quad S(r+\varepsilon / 2)
$$

as soon as $\ell$ is large enough, i.e. $(1+\varepsilon) / \ell \leqslant \varepsilon^{2} / 2$. Let $\chi_{\varepsilon}$ be a cut-off function equal to 1 in $B(r+\varepsilon / 2)$ with support in $B(r+\varepsilon)$. Then

$$
\begin{aligned}
\int_{\bar{B}(r)} T_{k} \wedge\left(d d^{c} \varphi\right)^{p} & \leqslant \int_{B(r+\varepsilon / 2)} T_{k} \wedge\left(d d^{c} \psi_{\ell}\right)^{p} \\
& =(1+\varepsilon)^{p} \int_{B(r+\varepsilon / 2)} T_{k} \wedge\left(d d^{c} \varphi_{\ell}\right)^{p} \\
& \leqslant(1+\varepsilon)^{p} \int_{B(r+\varepsilon)} \chi_{\varepsilon} T_{k} \wedge\left(d d^{c} \varphi_{\ell}\right)^{p} .
\end{aligned}
$$

As $\chi_{\varepsilon}\left(d d^{c} \varphi_{\ell}\right)^{p}$ is smooth with compact support and as $T_{k}$ converges weakly to $T$, we infer

$$
\limsup _{k \rightarrow+\infty} \int_{\bar{B}(r)} T_{k} \wedge\left(d d^{c} \varphi\right)^{p} \leqslant(1+\varepsilon)^{p} \int_{B(r+\varepsilon)} \chi_{\varepsilon} T \wedge\left(d d^{c} \varphi_{\ell}\right)^{p}
$$

We then let $\ell$ tend to $+\infty$ and $\varepsilon$ tend to 0 to get the desired inequality. The first inequality is obtained in a similar way, we define $\psi_{\ell}$ so that $\psi_{\ell}=\varphi$ on $X \backslash B(r)$ and $\psi_{\ell}=\max \left\{(1-\varepsilon)\left(\varphi_{\ell}-1 / \ell\right)+r \varepsilon\right\}$ on $\bar{B}(r)$, and we take $\chi_{\varepsilon}=1$ on $B(r-\varepsilon)$ with Supp $\chi_{\varepsilon} \subset B(r-\varepsilon / 2)$. Then for $\ell$ large

$$
\begin{aligned}
\int_{B(r)} T_{k} \wedge\left(d d^{c} \varphi\right)^{p} & \geqslant \int_{B(r-\varepsilon / 2)} T_{k} \wedge\left(d d^{c} \psi_{\ell}\right)^{p} \\
& \geqslant(1-\varepsilon)^{p} \int_{B(r-\varepsilon / 2)} \chi_{\varepsilon} T_{k} \wedge\left(d d^{c} \varphi_{\ell}\right)^{p} .
\end{aligned}
$$

(5.13) Proposition. Let $\varphi_{k}$ be a (non necessarily monotone) sequence of continuous plurisubharmonic functions such that $e^{\varphi_{k}}$ converges uniformly to $e^{\varphi}$ on every compact subset of $X$. Suppose that $\{\varphi<R\} \cap \operatorname{Supp} T \subset \subset X$. Then for $r<R$ we have

$$
\limsup _{k \rightarrow+\infty} \int_{\left\{\varphi_{k} \leqslant r\right\} \cap\{\varphi<R\}} T \wedge\left(d d^{c} \varphi_{k}\right)^{p} \leqslant \int_{\{\varphi \leqslant r\}} T \wedge\left(d d^{c} \varphi\right)^{p} .
$$

In particular $\lim \sup _{k \rightarrow+\infty} \nu\left(T, \varphi_{k}\right) \leqslant \nu(T, \varphi)$.
When we take $\varphi_{k}(z)=\log \left|z-a_{k}\right|$ with $a_{k} \rightarrow a$, Prop. 5.13 implies the upper semicontinuity of $a \mapsto \nu(T, a)$ which was already noticed in the proof of Prop. 5.11.

Proof. Our assumption is equivalent to saying that $\max \left\{\varphi_{k}, t\right\}$ converges locally uniformly to $\max \{\varphi, t\}$ for every $t$. Then Cor. 3.6 shows that $T \wedge\left(d d^{c} \max \left\{\varphi_{k}, t\right\}\right)^{p}$ converges weakly to $T \wedge\left(d d^{c} \max \{\varphi, t\}\right)^{p}$. If $\chi_{\varepsilon}$ is a cut-off function equal to 1 on $\{\varphi \leqslant r+\varepsilon / 2\}$ with support in $\{\varphi<r+\varepsilon\}$, we get

$$
\lim _{k \rightarrow+\infty} \int_{X} \chi_{\varepsilon} T \wedge\left(d d^{c} \max \left\{\varphi_{k}, t\right\}\right)^{p}=\int_{X} \chi_{\varepsilon} T \wedge\left(d d^{c} \max \{\varphi, t\}\right)^{p}
$$

For $k$ large, we have $\left\{\varphi_{k} \leqslant r\right\} \cap\{\varphi<R\} \subset\{\varphi<r+\varepsilon / 2\}$, thus when $\varepsilon$ tends to 0 we infer

$$
\limsup _{k \rightarrow+\infty} \int_{\left\{\varphi_{k} \leqslant r\right\} \cap\{\varphi<R\}} T \wedge\left(d d^{c} \max \left\{\varphi_{k}, t\right\}\right)^{p} \leqslant \int_{\{\varphi \leqslant r\}} T \wedge\left(d d^{c} \max \{\varphi, t\}\right)^{p}
$$

When we choose $t<r$, this is equivalent to the first inequality in statement (5.13).

## § 6. The Jensen-Lelong Formula

We assume in this section that $X$ is Stein, that $\varphi$ is semi-exhaustive on $X$ and that $B(R) \subset \subset$. We set for simplicity $\varphi_{\geqslant r}=\max \{\varphi, r\}$. For every $\left.r \in\right]-\infty, R[$, the measures $d d^{c}\left(\varphi_{\geqslant r}\right)^{n}$ are well defined. By Cor. 3.6, the map $r \longmapsto\left(d d^{c} \varphi_{\geqslant r}\right)^{n}$ is continuous on $]-\infty, R\left[\right.$ with respect to the weak topology. As $\left(d d^{c} \varphi_{\geqslant r}\right)^{n}=\left(d d^{c} \varphi\right)^{n}$ on $X \backslash \bar{B}(r)$ and as $\varphi_{\geqslant r} \equiv r,\left(d d^{c} \varphi_{\geqslant r}\right)^{n}=0$ on $B(r)$, the left continuity implies $\left(d d^{c} \varphi_{\geqslant r}\right)^{n} \geqslant$ $\mathbb{1}_{X \backslash B(r)}\left(d d^{c} \varphi\right)^{n}$. Here $\mathbb{1}_{A}$ denotes the characteristic function of any subset $A \subset X$. According to the definition introduced in [Demailly 1985a], the collection of MongeAmpère measures associated with $\varphi$ is the family of positive measures $\mu_{r}$ such that

$$
\begin{equation*}
\left.\mu_{r}=\left(d d^{c} \varphi_{\geqslant r}\right)^{n}-\mathbb{1}_{X \backslash B(r)}\left(d d^{c} \varphi\right)^{n}, \quad r \in\right]-\infty, R[. \tag{6.1}
\end{equation*}
$$

The measure $\mu_{r}$ is supported on $S(r)$ and $r \longmapsto \mu_{r}$ is weakly continuous on the left by the bounded convergence theorem. Stokes' formula shows that $\int_{B(s)}\left(d d^{c} \varphi \geqslant r\right)^{n}-\left(d d^{c} \varphi\right)^{n}=0$ for $s>r$, hence the total mass $\mu_{r}(S(r))=\mu_{r}(B(s))$ is equal to the difference between the masses of $\left(d d^{c} \varphi\right)^{n}$ and $\mathbb{1}_{X \backslash B(r)}\left(d d^{c} \varphi\right)^{n}$ over $B(s)$, i.e.

$$
\begin{equation*}
\mu_{r}(S(r))=\int_{B(r)}\left(d d^{c} \varphi\right)^{n} \tag{6.2}
\end{equation*}
$$

(6.3) Example. When $\left(d d^{c} \varphi\right)^{n}=0$ on $X \backslash \varphi^{-1}(-\infty)$, formula (6.1) can be simplified into $\mu_{r}=\left(d d^{c} \varphi_{\geqslant r}\right)^{n}$. This is so for $\varphi(z)=\log |z|$. In this case, the invariance of $\varphi$ under unitary transformations implies that $\mu_{r}$ is also invariant. As the total mass of $\mu_{r}$ is equal to 1 by 5.10 and (6.2), we see that $\mu_{r}$ is the invariant measure of mass 1 on the euclidean sphere of radius $e^{r}$.
(6.4) Proposition. Assume that $\varphi$ is smooth near $S(r)$ and that $d \varphi \neq 0$ on $S(r)$, i.e. $r$ is a non critical value. Then $S(r)=\partial B(r)$ is a smooth oriented real hypersurface and the measure $\mu_{r}$ is given by the $(2 n-1)$-volume form $\left(d d^{c} \varphi\right)^{n-1} \wedge d^{c} \varphi_{\lceil S(r)}$.

Proof. Write $\max \{t, r\}=\lim _{k \rightarrow+\infty} \chi_{k}(t)$ where $\chi$ is a decreasing sequence of smooth convex functions with $\chi_{k}(t)=r$ for $t \leqslant r-1 / k, \chi_{k}(t)=t$ for $t \geqslant r+1 / k$. Theorem 3.6 shows that $\left(d d^{c} \chi_{k} \circ \varphi\right)^{n}$ converges weakly to $\left(d d^{c} \varphi_{\geqslant r}\right)^{n}$. Let $h$ be a smooth function $h$ with compact support near $S(r)$. Let us apply Stokes' theorem with $S(r)$ considered as the boundary of $X \backslash B(r)$ :

$$
\begin{aligned}
\int_{X} h\left(d d^{c} \varphi_{\geqslant r}\right)^{n} & =\lim _{k \rightarrow+\infty} \int_{X} h\left(d d^{c} \chi_{k} \circ \varphi\right)^{n} \\
& =\lim _{k \rightarrow+\infty} \int_{X}-d h \wedge\left(d d^{c} \chi_{k} \circ \varphi\right)^{n-1} \wedge d^{c}\left(\chi_{k} \circ \varphi\right)
\end{aligned}
$$

$$
\begin{aligned}
& =\lim _{k \rightarrow+\infty} \int_{X}-\chi_{k}^{\prime}(\varphi)^{n} d h \wedge\left(d d^{c} \varphi\right)^{n-1} \wedge d^{c} \varphi \\
& =\int_{X \backslash B(r)}-d h \wedge\left(d d^{c} \varphi\right)^{n-1} \wedge d^{c} \varphi \\
& =\int_{S(r)} h\left(d d^{c} \varphi\right)^{n-1} \wedge d^{c} \varphi+\int_{X \backslash B(r)} h\left(d d^{c} \varphi\right)^{n-1} \wedge d d^{c} \varphi .
\end{aligned}
$$

Near $S(r)$ we thus have an equality of measures

$$
\left(d d^{c} \varphi_{\geqslant r}\right)^{n}=\left(d d^{c} \varphi\right)^{n-1} \wedge d^{c} \varphi_{\mid S(r)}+\mathbb{1}_{X \backslash B(r)}\left(d d^{c} \varphi\right)^{n} .
$$

(6.5) Jensen-Lelong formula. Let $V$ be any plurisubharmonic function on $X$. Then $V$ is $\mu_{r}$-integrable for every $\left.r \in\right]-\infty, R[$ and

$$
\mu_{r}(V)-\int_{B(r)} V\left(d d^{c} \varphi\right)^{n}=\int_{-\infty}^{r} \nu\left(d d^{c} V, \varphi, t\right) d t
$$

Proof. Proposition 3.11 shows that $V$ is integrable with respect to the measure $\left(d d^{c} \varphi_{\geqslant r}\right)^{n}$, hence $V$ is $\mu_{r}$-integrable. By definition

$$
\nu\left(d d^{c} V, \varphi, t\right)=\int_{\varphi(z)<t} d d^{c} V \wedge\left(d d^{c} \varphi\right)^{n-1}
$$

and the Fubini theorem gives

$$
\begin{align*}
\int_{-\infty}^{r} \nu\left(d d^{c} V, \varphi, t\right) d t & =\iint_{\varphi(z)<t<r} d d^{c} V(z) \wedge\left(d d^{c} \varphi(z)\right)^{n-1} d t \\
& =\int_{B(r)}(r-\varphi) d d^{c} V \wedge\left(d d^{c} \varphi\right)^{n-1} \tag{6.6}
\end{align*}
$$

We first show that Formula 6.5 is true when $\varphi$ and $V$ are smooth. As both members of the formula are left continuous with respect to $r$ and as almost all values of $\varphi$ are non critical by Sard's theorem, we may assume $r$ non critical. Formula 3.1 applied with $f=(r-\varphi)\left(d d^{c} \varphi\right)^{n-1}$ and $g=V$ shows that integral (6.6) is equal to

$$
\int_{S(r)} V\left(d d^{c} \varphi\right)^{n-1} \wedge d^{c} \varphi-\int_{B(r)} V\left(d d^{c} \varphi\right)^{n}=\mu_{r}(V)-\int_{B(r)} V\left(d d^{c} \varphi\right)^{n}
$$

Formula 6.5 is thus proved when $\varphi$ and $V$ are smooth. If $V$ is smooth and $\varphi$ merely continuous and finite, one can write $\varphi=\lim \varphi_{k}$ where $\varphi_{k}$ is a decreasing sequence of smooth plurisubharmonic functions (because $X$ is Stein). Then $d d^{c} V \wedge\left(d d^{c} \varphi_{k}\right)^{n-1}$ converges weakly to $d d^{c} V \wedge\left(d d^{c} \varphi\right)^{n-1}$ and (6.6) converges, since $\mathbb{1}_{B(r)}(r-\varphi)$ is continuous with compact support on $X$. The left hand side of Formula 6.5 also converges because the definition of $\mu_{r}$ implies

$$
\mu_{k, r}(V)-\int_{\varphi_{k}<r} V\left(d d^{c} \varphi_{k}\right)^{n}=\int_{X} V\left(\left(d d^{c} \varphi_{k, \geqslant r}\right)^{n}-\left(d d^{c} \varphi_{k}\right)^{n}\right)
$$

and we can apply again weak convergence on a neighborhood of $\bar{B}(r)$. If $\varphi$ takes $-\infty$ values, replace $\varphi$ by $\varphi \geqslant-k$ where $k \rightarrow+\infty$. Then $\mu_{r}(V)$ is unchanged, $\int_{B(r)} V\left(d d^{c} \varphi_{\geqslant-k}\right)^{n}$ converges to $\int_{B(r)} V\left(d d^{c} \varphi\right)^{n}$ and the right hand side of Formula 6.5 is replaced by $\int_{-k}^{r} \nu\left(d d^{c} V, \varphi, t\right) d t$. Finally, for $V$ arbitrary, write $V=\lim \downarrow V_{k}$ with a sequence of smooth functions $V_{k}$. Then $d d^{c} V_{k} \wedge\left(d d^{c} \varphi\right)^{n-1}$ converges weakly to $d d^{c} V \wedge\left(d d^{c} \varphi\right)^{n-1}$ by Prop. 4.4, so the integral (6.6) converges to the expected limit and the same is true for the left hand side of 6.5 by the monotone convergence theorem.

For $r<r_{0}<R$, the Jensen-Lelong formula implies

$$
\begin{equation*}
\mu_{r}(V)-\mu_{r_{0}}(V)+\int_{B\left(r_{0}\right) \backslash B(r)} V\left(d d^{c} \varphi\right)^{n}=\int_{r_{0}}^{r} \nu\left(d d^{c} V, \varphi, t\right) d t \tag{6.7}
\end{equation*}
$$

(6.8) Corollary. Assume that $\left(d d^{c} \varphi\right)^{n}=0$ on $X \backslash S(-\infty)$. Then $r \mapsto \mu_{r}(V)$ is a convex increasing function of $r$ and the lelong number $\nu\left(d d^{c} V, \varphi\right)$ is given by

$$
\nu\left(d d^{c} V, \varphi\right)=\lim _{r \rightarrow-\infty} \frac{\mu_{r}(V)}{r} .
$$

Proof. By (6.7) we have

$$
\mu_{r}(V)=\mu_{r_{0}}(V)+\int_{r_{0}}^{r} \nu\left(d d^{c} V, \varphi, t\right) d t
$$

As $\nu\left(d d^{c} V, \varphi, t\right)$ is increasing and nonnegative, it follows that $r \longmapsto \mu_{r}(V)$ is convex and increasing. The formula for $\nu\left(d d^{c} V, \varphi\right)=\lim _{t \rightarrow-\infty} \nu\left(d d^{c} V, \varphi, t\right)$ is then obvious.
(6.9) Example. Let $X$ be an open subset of $\mathbb{C}^{n}$ equipped with the semi-exhaustive function $\varphi(z)=\log |z-a|, a \in X$. Then $\left(d d^{c} \varphi\right)^{n}=\delta_{a}$ and the Jensen-Lelong formula becomes

$$
\mu_{r}(V)=V(a)+\int_{-\infty}^{r} \nu\left(d d^{c} V, \varphi, t\right) d t
$$

As $\mu_{r}$ is the mean value measure on the sphere $S\left(a, e^{r}\right)$, we make the change of variables $r \mapsto \log r, t \mapsto \log t$ and obtain the more familiar formula

$$
\begin{equation*}
\mu(V, S(a, r))=V(a)+\int_{0}^{r} \nu\left(d d^{c} V, a, t\right) \frac{d t}{t} \tag{6.9a}
\end{equation*}
$$

where $\nu\left(d d^{c} V, a, t\right)=\nu\left(d d^{c} V, \varphi, \log t\right)$ is given by (5.7):

$$
\begin{equation*}
\nu\left(d d^{c} V, a, t\right)=\frac{1}{\pi^{n-1} t^{2 n-2} /(n-1)!} \int_{B(a, t)} \frac{1}{2 \pi} \Delta V \tag{6.9b}
\end{equation*}
$$

In this setting, Cor. 6.8 implies

$$
\begin{equation*}
\nu\left(d d^{c} V, a\right)=\lim _{r \rightarrow 0} \frac{\mu(V, S(a, r))}{\log r}=\lim _{r \rightarrow 0} \frac{\sup _{S(a, r)} V}{\log r} \tag{6.9c}
\end{equation*}
$$

To prove the last equality, we may assume $V \leqslant 0$ after subtraction of a constant. Inequality $\geqslant$ follows from the obvious estimate $\mu(V, S(a, r)) \leqslant \sup _{S(a, r)} V$, while inequality $\leqslant$ follows from the standard Harnack estimate

$$
\begin{equation*}
\sup _{S(a, \varepsilon r)} V \leqslant \frac{1-\varepsilon}{(1+\varepsilon)^{2 n-1}} \mu(V, S(a, r)) \tag{6.9d}
\end{equation*}
$$

when $\varepsilon$ is small (this estimate follows easily from the Green-Riesz representation formula 1.4.6 and 1.4.7). As $\sup _{S(a, r)} V=\sup _{B(a, r)} V$, Formula (6.9c) can also be rewritten $\nu\left(d d^{c} V, a\right)=\liminf _{z \rightarrow a} V(z) / \log |z-a|$. Since $\sup _{S(a, r)} V$ is a convex (increasing) function of $\log r$, we infer that

$$
\begin{equation*}
V(z) \leqslant \gamma \log |z-a|+\mathrm{O}(1) \tag{6.9e}
\end{equation*}
$$

with $\gamma=\nu\left(d d^{c} V, a\right)$, and $\nu\left(d d^{c} V, a\right)$ is the largest constant $\gamma$ which satisfies this inequality. Thus $\nu\left(d d^{c} V, a\right)=\gamma$ is equivalent to $V$ having a logarithmic pole of coefficient $\gamma$.
§ 6.10. Special case. Take in particular $V=\log |f|$ where $f$ is a holomorphic function on $X$. The Lelong-Poincaré formula shows that $d d^{c} \log |f|$ is equal to the zero divisor $\left[Z_{f}\right]=\sum m_{j}\left[H_{j}\right]$, where $H_{j}$ are the irreducible components of $f^{-1}(0)$ and $m_{j}$ is the multiplicity of $f$ on $H_{j}$. The trace $\frac{1}{2 \pi} \Delta \log |f|$ is then the euclidean area measure of $Z_{f}$ (with corresponding multiplicities $m_{j}$ ). By Formula ( 6.9 c ), we see that the Lelong number $\nu\left(\left[Z_{f}\right], a\right)$ is equal to the vanishing order $\operatorname{ord}_{a}(f)$, that is, the smallest integer $m$ such that $D^{\alpha} f(a) \neq 0$ for some multiindex $\alpha$ with $|\alpha|=m$. In dimension $n=1$, we have $\frac{1}{2 \pi} \Delta \log f=\sum m_{j} \delta_{a_{j}}$. Then (6.9a) is the usual Jensen formula

$$
\mu(\log |f|, S(0, r))-\log |f(0)|=\int_{0}^{r} \nu(t) \frac{d t}{t}=\sum m_{j} \log \frac{r}{\left|a_{j}\right|}
$$

where $\nu(t)$ is the number of zeros $a_{j}$ in the disk $D(0, t)$, counted with multiplicities $m_{j}$.
(6.11) Example. Take $\varphi(z)=\log \max \left|z_{j}\right|^{\lambda_{j}}$ where $\lambda_{j}>0$. Then $B(r)$ is the polydisk of radii $\left(e^{r / \lambda_{1}}, \ldots, e^{r / \lambda_{n}}\right)$. If some coordinate $z_{j}$ is non zero, say $z_{1}$, we can write $\varphi(z)$ as $\lambda_{1} \log \left|z_{1}\right|$ plus some function depending only on the $(n-1)$ variables $z_{j} / z_{1}^{\lambda_{1} / \lambda_{j}}$. Hence $\left(d d^{c} \varphi\right)^{n}=0$ on $\mathbb{C}^{n} \backslash\{0\}$. It will be shown later that

$$
\begin{equation*}
\left(d d^{c} \varphi\right)^{n}=\lambda_{1} \ldots \lambda_{n} \delta_{0} \tag{6.11a}
\end{equation*}
$$

We now determine the measures $\mu_{r}$. At any point $z$ where not all terms $\left|z_{j}\right|^{\lambda_{j}}$ are equal, the smallest one can be omitted without changing $\varphi$ in a neighborhood of $z$. Thus $\varphi$ depends only on $(n-1)$-variables and $\left(d d^{c} \varphi_{\geqslant r}\right)^{n}=0, \mu_{r}=0$ near $z$. It follows that $\mu_{r}$ is supported by the distinguished boundary $\left|z_{j}\right|=e^{r / \lambda_{j}}$ of the polydisk $B(r)$. As $\varphi$ is invariant by all rotations $z_{j} \longmapsto e^{\mathrm{i} \theta_{j}} z_{j}$, the measure $\mu_{r}$ is also invariant and we see that $\mu_{r}$ is a constant multiple of $d \theta_{1} \ldots d \theta_{n}$. By formula (6.2) and (6.11 a) we get

$$
\begin{equation*}
\mu_{r}=\lambda_{1} \ldots \lambda_{n}(2 \pi)^{-n} d \theta_{1} \ldots d \theta_{n} \tag{6.11b}
\end{equation*}
$$

In particular, the Lelong number $\nu\left(d d^{c} V, \varphi\right)$ is given by

$$
\nu\left(d d^{c} V, \varphi\right)=\lim _{r \rightarrow-\infty} \frac{\lambda_{1} \ldots \lambda_{n}}{r} \int_{\theta_{j} \in[0,2 \pi]} V\left(e^{r / \lambda_{1}+\mathrm{i} \theta_{1}}, \ldots, e^{r / \lambda_{n}+\mathrm{i} \theta_{n}}\right) \frac{d \theta_{1} \ldots d \theta_{n}}{(2 \pi)^{n}}
$$

These numbers have been introduced and studied by [Kiselman 1986]. We call them directional Lelong numbers with coefficients $\left(\lambda_{1}, \ldots, \lambda_{n}\right)$. For an arbitrary current $T$, we define

$$
\begin{equation*}
\nu(T, x, \lambda)=\nu\left(T, \log \max \left|z_{j}-x_{j}\right|^{\lambda_{j}}\right) \tag{6.11c}
\end{equation*}
$$

The above formula for $\nu\left(d d^{c} V, \varphi\right)$ combined with the analogue of Harnack's inequality (6.9 d) for polydisks gives

$$
\begin{align*}
\nu\left(d d^{c} V, x, \lambda\right) & =\lim _{r \rightarrow 0} \frac{\lambda_{1} \ldots \lambda_{n}}{\log r} \int V\left(r^{1 / \lambda_{1}} e^{\mathrm{i} \theta_{1}}, \ldots, r^{1 / \lambda_{n}} e^{\mathrm{i} \theta_{n}}\right) \frac{d \theta_{1} \ldots d \theta_{n}}{(2 \pi)^{n}} \\
& =\lim _{r \rightarrow 0} \frac{\lambda_{1} \ldots \lambda_{n}}{\log r} \sup _{\theta_{1}, \ldots, \theta_{n}} V\left(r^{1 / \lambda_{1}} e^{\mathrm{i} \theta_{1}}, \ldots, r^{1 / \lambda_{n}} e^{\mathrm{i} \theta_{n}}\right) . \tag{6.11~d}
\end{align*}
$$

## § 7. Comparison Theorems for Lelong Numbers

Let $T$ be a closed positive current of bidimension $(p, p)$ on a Stein manifold $X$ equipped with a semi-exhaustive plurisubharmonic weight $\varphi$. We first show that the Lelong numbers $\nu(T, \varphi)$ only depend on the asymptotic behaviour of $\varphi$ near the polar set $S(-\infty)$. In a precise way:
(7.1) First comparison theorem. Let $\varphi, \psi: X \longrightarrow[-\infty,+\infty[$ be continuous plurisubharmonic functions. We assume that $\varphi, \psi$ are semi-exhaustive on Supp $T$ and that

$$
\ell:=\lim \sup \frac{\psi(x)}{\varphi(x)}<+\infty \quad \text { as } \quad x \in \operatorname{Supp} T \text { and } \varphi(x) \rightarrow-\infty
$$

Then $\nu(T, \psi) \leqslant \ell^{p} \nu(T, \varphi)$, and the equality holds if $\ell=\lim \psi / \varphi$.
Proof. Definition 6.4 shows immediately that $\nu(T, \lambda \varphi)=\lambda^{p} \nu(T, \varphi)$ for every scalar $\lambda>0$. It is thus sufficient to verify the inequality $\nu(T, \psi) \leqslant \nu(T, \varphi)$ under the hypothesis $\lim \sup \psi / \varphi<1$. For all $c>0$, consider the plurisubharmonic function

$$
u_{c}=\max (\psi-c, \varphi) .
$$

Let $R_{\varphi}$ and $R_{\psi}$ be such that $B_{\varphi}\left(R_{\varphi}\right) \cap \operatorname{Supp} T$ and $B_{\psi}\left(R_{\psi}\right) \cap \operatorname{Supp} T$ be relatively compact in $X$. Let $r<R_{\varphi}$ and $a<r$ be fixed. For $c>0$ large enough, we have $u_{c}=\varphi$ on $\varphi^{-1}([a, r])$ and Stokes' formula gives

$$
\nu(T, \varphi, r)=\nu\left(T, u_{c}, r\right) \geqslant \nu\left(T, u_{c}\right)
$$

The hypothesis $\lim \sup \psi / \varphi<1$ implies on the other hand that there exists $t_{0}<0$ such that $u_{c}=\psi-c$ on $\left\{u_{c}<t_{0}\right\} \cap$ Supp $T$. We infer

$$
\nu\left(T, u_{c}\right)=\nu(T, \psi-c)=\nu(T, \psi)
$$

hence $\nu(T, \psi) \leqslant \nu(T, \varphi)$. The equality case is obtained by reversing the roles of $\varphi$ and $\psi$ and observing that $\lim \varphi / \psi=1 / l$.

Assume in particular that $z^{k}=\left(z_{1}^{k}, \ldots, z_{n}^{k}\right), k=1,2$, are coordinate systems centered at a point $x \in X$ and let

$$
\varphi_{k}(z)=\log \left|z^{k}\right|=\log \left(\left|z_{1}^{k}\right|^{2}+\ldots+\left|z_{n}^{k}\right|^{2}\right)^{1 / 2}
$$

We have $\lim _{z \rightarrow x} \varphi_{2}(z) / \varphi_{1}(z)=1$, hence $\nu\left(T, \varphi_{1}\right)=\nu\left(T, \varphi_{2}\right)$ by Th. 7.1.
(7.2) Corollary. The usual Lelong numbers $\nu(T, x)$ are independent of the choice of local coordinates.

This result had been originally proved by [Siu 1974] with a much more delicate proof. Another interesting consequence is:
(7.3) Corollary. On an open subset of $\mathbb{C}^{n}$, the Lelong numbers and Kiselman numbers are related by

$$
\nu(T, x)=\nu(T, x,(1, \ldots, 1))
$$

Proof. By definition, the Lelong number $\nu(T, x)$ is associated with the weight $\varphi(z)=$ $\log |z-x|$ and the Kiselman number $\nu(T, x,(1, \ldots, 1))$ to the weight $\psi(z)=\log \max \mid z_{j}-$ $x_{j} \mid$. It is clear that $\lim _{z \rightarrow x} \psi(z) / \varphi(z)=1$, whence the conclusion.

Another consequence of Th. 7.1 is that $\nu(T, x, \lambda)$ is an increasing function of each variable $\lambda_{j}$. Moreover, if $\lambda_{1} \leqslant \ldots \leqslant \lambda_{n}$, we get the inequalities

$$
\lambda_{1}^{p} \nu(T, x) \leqslant \nu(T, x, \lambda) \leqslant \lambda_{n}^{p} \nu(T, x) .
$$

These inequalities will be improved in section 7 (see Cor. 9.16). For the moment, we just prove the following special case.
(7.4) Corollary. For all $\lambda_{1}, \ldots, \lambda_{n}>0$ we have

$$
\left(d d^{c} \log \max _{1 \leqslant j \leqslant n}\left|z_{j}\right|^{\lambda_{j}}\right)^{n}=\left(d d^{c} \log \sum_{1 \leqslant j \leqslant n}\left|z_{j}\right|^{\lambda_{j}}\right)^{n}=\lambda_{1} \ldots \lambda_{n} \delta_{0} .
$$

Proof. In fact, our measures vanish on $\mathbb{C}^{n} \backslash\{0\}$ by the arguments explained in example 6.11. Hence they are equal to $c \delta_{0}$ for some constant $c \geqslant 0$ which is simply the Lelong number of the bidimension $(n, n)$-current $T=[X]=1$ with the corresponding weight. The comparison theorem shows that the first equality holds and that

$$
\left(d d^{c} \log \sum_{1 \leqslant j \leqslant n}\left|z_{j}\right|^{\lambda_{j}}\right)^{n}=\ell^{-n}\left(d d^{c} \log \sum_{1 \leqslant j \leqslant n}\left|z_{j}\right|^{\mid \lambda_{j}}\right)^{n}
$$

for all $\ell>0$. By taking $\ell$ large and approximating $\ell \lambda_{j}$ with $2\left[\ell \lambda_{j} / 2\right]$, we may assume that $\lambda_{j}=2 s_{j}$ is an even integer. Then formula (5.6) gives

$$
\begin{aligned}
& \int_{\sum\left|z_{j}\right|^{2 s_{j}}<r^{2}}\left(d d^{c} \log \sum\left|z_{j}\right|^{2 s_{j}}\right)^{n}=r^{-2 n} \int_{\sum\left|z_{j}\right|^{2 s_{j}}<r^{2}}\left(d d^{c} \sum\left|z_{j}\right|^{2 s_{j}}\right)^{n} \\
& \quad=s_{1} \ldots s_{n} r^{-2 n} \int_{\sum\left|w_{j}\right|^{2}<r^{2}} 2^{n}\left(\frac{\mathrm{i}}{2 \pi} d^{\prime} d^{\prime \prime}|w|^{2}\right)^{n}=\lambda_{1} \ldots \lambda_{n}
\end{aligned}
$$

by using the $s_{1} \ldots s_{n}$-sheeted change of variables $w_{j}=z_{j}^{s_{j}}$.
Now, we assume that $T=[A]$ is the current of integration over an analytic set $A \subset X$ of pure dimension $p$. The above comparison theorem will enable us to give a simple proof of P. Thie's main result [Thie 1967]: the Lelong number $\nu([A], x)$ can be interpreted as the multiplicity of the analytic set $A$ at point $x$. Our starting point is the following consequence of Th. II.3.19 applied simultaneously to all irreducible components of $(A, x)$.
(7.5) Lemma. For a generic choice of local coordinates $z^{\prime}=\left(z_{1}, \ldots, z_{p}\right)$ and $z^{\prime \prime}=\left(z_{p+1}, \ldots, z_{n}\right)$ on $(X, x)$, the germ $(A, x)$ is contained in a cone $\left|z^{\prime \prime}\right| \leqslant C\left|z^{\prime}\right|$. If $B^{\prime} \subset \mathbb{C}^{p}$ is a ball of center 0 and radius $r^{\prime}$ small, and $B^{\prime \prime} \subset \mathbb{C}^{n-p}$ is the ball of center 0 and radius $r^{\prime \prime}=C r^{\prime}$, then the projection

$$
\operatorname{pr}: A \cap\left(B^{\prime} \times B^{\prime \prime}\right) \longrightarrow B^{\prime}
$$

is a ramified covering with finite sheet number $m$.

We use these properties to compute the Lelong number of $[A]$ at point $x$. When $z \in A$ tends to $x$, the functions

$$
\varphi(z)=\log |z|=\log \left(\left|z^{\prime}\right|^{2}+\left|z^{\prime \prime}\right|^{2}\right)^{1 / 2}, \quad \psi(z)=\log \left|z^{\prime}\right| .
$$

are equivalent. As $\varphi, \psi$ are semi-exhaustive on $A$, Th. 7.1 implies

$$
\nu([A], x)=\nu([A], \varphi)=\nu([A], \psi)
$$

Let us apply formula (5.6) to $\psi$ : for every $t<r^{\prime}$ we get

$$
\begin{aligned}
\nu([A], \psi, \log t) & =t^{-2 p} \int_{\{\psi<\log t\}}[A] \wedge\left(\frac{1}{2} d d^{c} e^{2 \psi}\right)^{p} \\
& =t^{-2 p} \int_{A \cap\left\{\left|z^{\prime}\right|<t\right\}}\left(\frac{1}{2} \operatorname{pr}^{\star} d d^{c}\left|z^{\prime}\right|^{2}\right)^{p} \\
& =m t^{-2 p} \int_{\mathbb{C}^{p} \cap\left\{\left|z^{\prime}\right|<t\right\}}\left(\frac{1}{2} d d^{c}\left|z^{\prime}\right|^{2}\right)^{p}=m,
\end{aligned}
$$

hence $\nu([A], \psi)=m$. Here, we have used the fact that pr is an étale covering with $m$ sheets over the complement of the ramification locus $S \subset B^{\prime}$, and the fact that $S$ is of zero Lebesgue measure in $B^{\prime}$. We have thus obtained simultaneously the following two results:
(7.6) Theorem and Definition. Let $A$ be an analytic set of dimension $p$ in a complex manifold $X$ of dimension $n$. For a generic choice of local coordinates $z^{\prime}=\left(z_{1}, \ldots, z_{p}\right)$, $z^{\prime \prime}=\left(z_{p+1}, \ldots, z_{n}\right)$ near a point $x \in A$ such that the germ $(A, x)$ is contained in a cone $\left|z^{\prime \prime}\right| \leqslant C\left|z^{\prime}\right|$, the sheet number $m$ of the projection $(A, x) \rightarrow\left(\mathbb{C}^{p}, 0\right)$ onto the first $p$ coordinates is independent of the choice of $z^{\prime}, z^{\prime \prime}$. This number $m$ is called the multiplicity of $A$ at $x$.
(7.7) Theorem ([Thie 1967]). One has $\nu([A], x)=m$.

There is another interesting version of the comparison theorem which compares the Lelong numbers of two currents obtained as intersection products (in that case, we take the same weight for both).
(7.8) Second comparison theorem. Let $u_{1}, \ldots, u_{q}$ and $v_{1}, \ldots, v_{q}$ be plurisubharmonic functions such that each $q$-tuple satisfies the hypotheses of $T h .4 .5$ with respect to $T$. Suppose moreover that $u_{j}=-\infty$ on Supp $T \cap \varphi^{-1}(-\infty)$ and that

$$
\ell_{j}:=\lim \sup \frac{v_{j}(z)}{u_{j}(z)}<+\infty \quad \text { when } \quad z \in \operatorname{Supp} T \backslash u_{j}^{-1}(-\infty), \quad \varphi(z) \rightarrow-\infty .
$$

Then

$$
\nu\left(d d^{c} v_{1} \wedge \ldots \wedge d d^{c} v_{q} \wedge T, \varphi\right) \leqslant \ell_{1} \ldots \ell_{q} \nu\left(d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T, \varphi\right)
$$

Proof. By homogeneity in each factor $v_{j}$, it is enough to prove the inequality with constants $\ell_{j}=1$ under the hypothesis $\lim \sup v_{j} / u_{j}<1$. We set

$$
w_{j, c}=\max \left\{v_{j}-c, u_{j}\right\}
$$

Our assumption implies that $w_{j, c}$ coincides with $v_{j}-c$ on a neighborhood Supp $T \cap\{\varphi<$ $\left.r_{0}\right\}$ of Supp $T \cap\{\varphi<-\infty\}$, thus

$$
\nu\left(d d^{c} v_{1} \wedge \ldots \wedge d d^{c} v_{q} \wedge T, \varphi\right)=\nu\left(d d^{c} w_{1, c} \wedge \ldots \wedge d d^{c} w_{q, c} \wedge T, \varphi\right)
$$

for every $c$. Now, fix $r<R_{\varphi}$. Proposition 4.9 shows that the current $d d^{c} w_{1, c} \wedge \ldots \wedge$ $d d^{c} w_{q, c} \wedge T$ converges weakly to $d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T$ when $c$ tends to $+\infty$. By Prop. 5.12 we get

$$
\limsup _{c \rightarrow+\infty} \nu\left(d d^{c} w_{1, c} \wedge \ldots \wedge d d^{c} w_{q, c} \wedge T, \varphi\right) \leqslant \nu\left(d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T, \varphi\right)
$$

(7.9) Corollary. If $d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T$ is well defined, then at every point $x \in X$ we have

$$
\nu\left(d d^{c} u_{1} \wedge \ldots \wedge d d^{c} u_{q} \wedge T, x\right) \geqslant \nu\left(d d^{c} u_{1}, x\right) \ldots \nu\left(d d^{c} u_{q}, x\right) \nu(T, x)
$$

Proof. Apply (7.8) with $\varphi(z)=v_{1}(z)=\ldots=v_{q}(z)=\log |z-x|$ and observe that $\ell_{j}:=\limsup v_{j} / u_{j}=1 / \nu\left(d d^{c} u_{j}, x\right)$ (there is nothing to prove if $\nu\left(d d^{c} u_{j}, x\right)=0$ ).

Finally, we present an interesting stability property of Lelong numbers due to [Siu 1974]: almost all slices of a closed positive current $T$ along linear subspaces passing through a given point have the same Lelong number as $T$. Before giving a proof of this, we need a useful formula known as Crofton's formula.
(7.10) Lemma. Let $\alpha$ be a closed positive $(p, p)$-form on $\mathbb{C}^{n} \backslash\{0\}$ which is invariant under the unitary group $U(n)$. Then $\alpha$ has the form

$$
\alpha=\left(d d^{c} \chi(\log |z|)\right)^{p}
$$

where $\chi$ is a convex increasing function. Moreover $\alpha$ is invariant by homotheties if and only if $\chi$ is an affine function, i.e. $\alpha=\lambda\left(d d^{c} \log |z|\right)^{p}$.

Proof. A radial convolution $\alpha_{\varepsilon}(z)=\int_{\mathbb{R}} \rho(t / \varepsilon) \alpha\left(e^{t} z\right) d t$ produces a smooth form with the same properties as $\alpha$ and $\lim _{\varepsilon \rightarrow 0} \alpha_{\varepsilon}=\alpha$. Hence we can suppose that $\alpha$ is smooth on $\mathbb{C}^{n} \backslash\{0\}$. At a point $z=\left(0, \ldots, 0, z_{n}\right)$, the $(p, p)$-form $\alpha(z) \in \bigwedge^{p, p}\left(\mathbb{C}^{n}\right)^{\star}$ must be invariant by $U(n-1)$ acting on the first $(n-1)$ coordinates. We claim that the subspace of $U(n-1)$-invariants in $\bigwedge^{p, p}\left(\mathbb{C}^{n}\right)^{\star}$ is generated by $\left(d d^{c}|z|^{2}\right)^{p}$ and $\left(d d^{c}|z|^{2}\right)^{p-1} \wedge \mathrm{i} d z_{n} \wedge d \bar{z}_{n}$. In fact, a form $\beta=\sum \beta_{I, J} d z_{I} \wedge d \bar{z}_{J}$ is invariant by $U(1)^{n-1} \subset U(n-1)$ if and only if $\beta_{I, J}=0$ for $I \neq J$, and invariant by the permutation group $\mathscr{S}_{n-1} \subset U(n-1)$ if and only if all coefficients $\beta_{I, I}$ (resp. $\beta_{J n, J n}$ ) with $I, J \subset\{1, \ldots, n-1\}$ are equal. Hence

$$
\beta=\lambda \sum_{|I|=p} d z_{I} \wedge d \bar{z}_{I}+\mu\left(\sum_{|J|=p-1} d z_{J} \wedge d \bar{z}_{J}\right) \wedge d z_{n} \wedge d \bar{z}_{n} .
$$

This proves our claim. As $d|z|^{2} \wedge d^{c}|z|^{2}=\frac{\mathrm{i}}{\pi}\left|z_{n}\right|^{2} d z_{n} \wedge d \bar{z}_{n}$ at $\left(0, \ldots, 0, z_{n}\right)$, we conclude that

$$
\alpha(z)=f(z)\left(d d^{c}|z|^{2}\right)^{p}+g(z)\left(d d^{c}|z|^{2}\right)^{p-1} \wedge d|z|^{2} \wedge d^{c}|z|^{2}
$$

for some smooth functions $f, g$ on $\mathbb{C}^{n} \backslash\{0\}$. The $U(n)$-invariance of $\alpha$ shows that $f$ and $g$ are radial functions. We may rewrite the last formula as

$$
\alpha(z)=u(\log |z|)\left(d d^{c} \log |z|\right)^{p}+v(\log |z|)\left(d d^{c} \log |z|\right)^{p-1} \wedge d \log |z| \wedge d^{c} \log |z|
$$

Here $\left(d d^{c} \log |z|\right)^{p}$ is a positive $(p, p)$-form coming from $\mathbb{P}^{n-1}$, hence it has zero contraction in the radial direction, while the contraction of the form $\left(d d^{c} \log |z|\right)^{p-1} \wedge d \log |z| \wedge$ $d^{c} \log |z|$ by the radial vector field is $\left(d d^{c} \log |z|\right)^{p-1}$. This shows easily that $\alpha(z) \geqslant 0$ if and only if $u, v \geqslant 0$. Next, the closedness condition $d \alpha=0$ gives $u^{\prime}-v=0$. Thus $u$ is increasing and we define a convex increasing function $\chi$ by $\chi^{\prime}=u^{1 / p}$. Then $v=u^{\prime}=$ $p \chi^{\prime p-1} \chi^{\prime \prime}$ and

$$
\alpha(z)=\left(d d^{c} \chi(\log |z|)\right)^{p} .
$$

If $\alpha$ is invariant by homotheties, the functions $u$ and $v$ must be constant, thus $v=0$ and $\alpha=\lambda\left(d d^{c} \log |z|\right)^{p}$.
(7.11) Corollary (Crofton's formula). Let $d v$ be the unique $U(n)$-invariant measure of mass 1 on the Grassmannian $G(p, n)$ of $p$-dimensional subspaces in $\mathbb{C}^{n}$. Then

$$
\int_{S \in G(p, n)}[S] d v(S)=\left(d d^{c} \log |z|\right)^{n-p}
$$

Proof. The left hand integral is a closed positive bidegree ( $n-p, n-p$ ) current which is invariant by $U(n)$ and by homotheties. By Lemma 7.10, this current must coincide with the form $\lambda\left(d d^{c} \log |z|\right)^{n-p}$ for some $\lambda \geqslant 0$. The coefficient $\lambda$ is the Lelong number at 0 . As $\nu([S], 0)=1$ for every $S$, we get $\lambda=\int_{G(p, n)} d v(S)=1$.

We now recall a few basic facts of slicing theory; see [Federer 1969] for details. Let $\sigma: M \rightarrow M^{\prime}$ be a submersion of smooth differentiable manifolds and let $\Theta$ be a locally flat current on $M$, that is a current which can be written locally as $\Theta=U+d V$ where
$U, V$ have locally integrable coefficients. It can be shown that every current $\Theta$ such that both $\Theta$ and $d \Theta$ have measure coefficients is locally flat; in particular, closed positive currents are locally flats. Then, for almost every $x^{\prime} \in M^{\prime}$, there is a well defined slice $\Theta_{x^{\prime}}$, which is the current on the fiber $\sigma^{-1}\left(x^{\prime}\right)$ defined by

$$
\Theta_{x^{\prime}}=U_{\left\lceil\sigma^{-1}\left(x^{\prime}\right)\right.}+d V_{\left\lceil\sigma^{-1}\left(x^{\prime}\right)\right.} .
$$

The restrictions of $U, V$ to the fibers exist for almost all $x^{\prime}$ by the Fubini theorem. It is easy to show by a regularization $\Theta_{\varepsilon}=\Theta \star \rho_{\varepsilon}$ that the slices of a closed positive current are again closed and positive: in fact $U_{\varepsilon, x^{\prime}}$ and $V_{\varepsilon, x^{\prime}}$ converge to $U_{x^{\prime}}$ and $V_{x^{\prime}}$ in $L_{\mathrm{loc}}^{1}$, thus $\Theta_{\varepsilon, x^{\prime}}$ converges weakly to $\Theta_{x^{\prime}}$ for almost every $x^{\prime}$. This kind of slicing can be referred to as parallel slicing (if we think of $\sigma$ as being a projection map). The kind of slicing we need (where the slices are taken over linear subspaces passing through a given point) is of a slightly different nature and is called concurrent slicing.

The possibility of concurrent slicing is proved as follows. Let $T$ be a closed positive current of bidimension $(p, p)$ in the ball $B(0, R) \subset \mathbb{C}^{n}$. Let

$$
Y=\left\{(x, S) \in \mathbb{C}^{n} \times G(q, n) ; x \in S\right\}
$$

be the total space of the tautological rank $q$ vector bundle over the Grassmannian $G(q, n)$, equipped with the obvious projections

$$
\sigma: Y \longrightarrow G(q, n), \quad \pi: Y \longrightarrow \mathbb{C}^{n}
$$

We set $Y_{R}=\pi^{-1}(B(0, R))$ and $Y_{R}^{\star}=\pi^{-1}(B(0, R) \backslash\{0\})$. The restriction $\pi_{0}$ of $\pi$ to $Y_{R}^{\star}$ is a submersion onto $B(0, R) \backslash\{0\}$, so we have a well defined pull-back $\pi_{0}^{\star} T$ over $Y_{R}^{\star}$. We would like to extend it as a pull-back $\pi^{\star} T$ over $Y_{R}$, so as to define slices $T_{\upharpoonright S}=\left(\pi^{\star} T\right)_{\mid \sigma^{-1}(S)}$; of course, these slices can be non zero only if the dimension of $S$ is at least equal to the degree of $T$, i.e. if $q \geqslant n-p$. We first claim that $\pi_{0}^{\star} T$ has locally finite mass near the zero section $\pi^{-1}(0)$ of $\sigma$. In fact let $\omega_{G}$ be a unitary invariant Kähler metric over $G(q, n)$ and let $\beta=d d^{c}|z|^{2}$ in $\mathbb{C}^{n}$. Then we get a Kähler metric on $Y$ defined by $\omega_{Y}=\sigma^{\star} \omega_{G}+\pi^{\star} \beta$. If $N=(q-1)(n-q)$ is the dimension of the fibers of $\pi$, the projection formula $\pi_{\star}\left(u \wedge \pi^{\star} v\right)=\left(\pi_{\star} u\right) \wedge v$ gives

$$
\pi_{\star} \omega_{Y}^{N+p}=\sum_{1 \leqslant k \leqslant p}\binom{N+p}{k} \beta^{k} \wedge \pi_{\star}\left(\sigma^{\star} \omega_{G}^{N+p-k}\right)
$$

Here $\pi_{\star}\left(\sigma^{\star} \omega_{G}^{N+p-k}\right)$ is a unitary and homothety invariant $(p-k, p-k)$ closed positive form on $\mathbb{C}^{n} \backslash\{0\}$, so $\pi_{\star}\left(\sigma^{\star} \omega_{G}^{N+p-k}\right)$ is proportional to $\left(d d^{c} \log |z|\right)^{p-k}$. With some constants $\lambda_{k}>0$, we thus get

$$
\begin{aligned}
\int_{Y_{r}^{\star}} \pi_{0}^{\star} T \wedge \omega_{Y}^{N+p} & =\sum_{0 \leqslant k \leqslant p} \lambda_{k} \int_{B(0, r) \backslash\{0\}} T \wedge \beta^{k} \wedge\left(d d^{c} \log |z|\right)^{p-k} \\
& =\sum_{0 \leqslant k \leqslant p} \lambda_{k} 2^{-(p-k)} r^{-2(p-k)} \int_{B(0, r) \backslash\{0\}} T \wedge \beta^{p}<+\infty
\end{aligned}
$$

The Skoda-El Mir theorem 2.3 shows that the trivial extension $\widetilde{\pi}_{0}^{\star} T$ of $\pi_{0}^{\star} T$ is a closed positive current on $Y_{R}$. Of course, the zero section $\pi^{-1}(0)$ might also carry some extra
mass of the desired current $\pi^{\star} T$. Since $\pi^{-1}(0)$ has codimension $q$, this extra mass cannot exist when $q>n-p=\operatorname{codim} \pi^{\star} T$ and we simply set $\pi^{\star} T=\widetilde{\pi}_{0}^{\star} T$. On the other hand, if $q=n-p$, we set

$$
\begin{equation*}
\pi^{\star} T:=\widetilde{\pi}_{0}^{\star} T+\nu(T, 0)\left[\pi^{-1}(0)\right] \tag{7.12}
\end{equation*}
$$

We can now apply parallel slicing with respect to $\sigma: Y_{R} \rightarrow G(q, n)$, which is a submersion: for almost all $S \in G(q, n)$, there is a well defined slice $T_{\upharpoonright S}=\left(\pi^{\star} T\right)_{\mid \sigma^{-1}(S)}$. These slices coincide with the usual restrictions of $T$ to $S$ if $T$ is smooth.
(7.13) Theorem ([Siu 1974]). For almost all $S \in G(q, n)$ with $q \geqslant n-p$, the slice $T_{\lceil S}$ satisfies $\nu\left(T_{\upharpoonright S}, 0\right)=\nu(T, 0)$.

Proof. If $q=n-p$, the slice $T_{\uparrow S}$ consists of some positive measure with support in $S \backslash\{0\}$ plus a Dirac measure $\nu(T, 0) \delta_{0}$ coming from the slice of $\nu(T, 0)\left[\pi^{-1}(0)\right]$. The equality $\nu\left(T_{\upharpoonright S}, 0\right)=\nu(T, 0)$ thus follows directly from (7.12).

In the general case $q>n-p$, it is clearly sufficient to prove the following two properties:
a) $\nu(T, 0, r)=\int_{S \in G(q, n)} \nu\left(T_{\upharpoonright S}, 0, r\right) d v(S) \quad$ for all $\left.r \in\right] 0, R[$;
b) $\quad \nu\left(T_{\lceil S}, 0\right) \geqslant \nu(T, 0) \quad$ for almost all $S$.

In fact, a) implies that $\nu(T, 0)$ is the average of all Lelong numbers $\nu\left(T_{\uparrow S}, 0\right)$ and the conjunction with b) implies that these numbers must be equal to $\nu(T, 0)$ for almost all $S$. In order to prove a) and b), we can suppose without loss of generality that $T$ is smooth on $B(0, R) \backslash\{0\}$. Otherwise, we perform a small convolution with respect to the action of $\mathrm{Gl}_{n}(\mathbb{C})$ on $\mathbb{C}^{n}$ :

$$
T_{\varepsilon}=\int_{g \in \mathrm{Gl}_{n}(\mathbb{C})} \rho_{\varepsilon}(g) g^{\star} T d v(g)
$$

where $\left(\rho_{\varepsilon}\right)$ is a regularizing family with support in an $\varepsilon$-neighborhood of the unit element of $\mathrm{Gl}_{n}(\mathbb{C})$. Then $T_{\varepsilon}$ is smooth in $B(0,(1-\varepsilon) R) \backslash\{0\}$ and converges weakly to $T$. Moreover, we have $\nu\left(T_{\varepsilon}, 0\right)=\nu(T, 0)$ by (7.2) and $\nu\left(T_{\uparrow S}, 0\right) \geqslant \lim \sup _{\varepsilon \rightarrow 0} \nu\left(T_{\varepsilon,\lceil S}, 0\right)$ by (5.12), thus a), b) are preserved in the limit. If $T$ is smooth on $B(0, R) \backslash\{0\}$, the slice $T_{\upharpoonright S}$ is defined for all $S$ and is simply the restriction of $T$ to $S \backslash\{0\}$ (carrying no mass at the origin).
a) Here we may even assume that $T$ is smooth at 0 by performing an ordinary convolution. As $T_{\upharpoonright S}$ has bidegree ( $n-p, n-p$ ), we have

$$
\nu\left(T_{\upharpoonright S}, 0, r\right)=\int_{S \cap B(0, r)} T \wedge \alpha_{S}^{q-(n-p)}=\int_{B(0, r)} T \wedge[S] \wedge \alpha_{S}^{p+q-n}
$$

where $\alpha_{S}=d d^{c} \log |w|$ and $w=\left(w_{1}, \ldots, w_{q}\right)$ are orthonormal coordinates on $S$. We simply have to check that

$$
\int_{S \in G(q, n)}[S] \wedge \alpha_{S}^{p+q-n} d v(S)=\left(d d^{c} \log |z|\right)^{p}
$$

However, both sides are unitary and homothety invariant ( $p, p$ )-forms with Lelong number 1 at the origin, so they must coincide by Lemma 7.11.
b) We prove the inequality when $S=\mathbb{C}^{q} \times\{0\}$. By the comparison theorem 7.1, for every $r>0$ and $\varepsilon>0$ we have

$$
\begin{align*}
& \int_{B(0, r)} T \wedge \gamma_{\varepsilon}^{p} \geqslant \nu(T, 0) \quad \text { where }  \tag{7.14}\\
& \gamma_{\varepsilon}=\frac{1}{2} d d^{c} \log \left(\varepsilon\left|z_{1}\right|^{2}+\ldots+\varepsilon\left|z_{q}\right|^{2}+\left|z_{q+1}\right|^{2}+\ldots+\left|z_{n}\right|^{2}\right)
\end{align*}
$$

We claim that the current $\gamma_{\varepsilon}^{p}$ converges weakly to

$$
[S] \wedge \alpha_{S}^{p+q-n}=[S] \wedge\left(\frac{1}{2} d d^{c} \log \left(\left|z_{1}\right|^{2}+\ldots+\left|z_{q}\right|^{2}\right)\right)^{p+q-n}
$$

as $\varepsilon$ tends to 0 . In fact, the Lelong number of $\gamma_{\varepsilon}^{p}$ at 0 is 1 , hence by homogeneity

$$
\int_{B(0, r)} \gamma_{\varepsilon}^{p} \wedge\left(d d^{c}|z|^{2}\right)^{n-p}=\left(2 r^{2}\right)^{n-p}
$$

for all $\varepsilon, r>0$. Therefore the family $\left(\gamma_{\varepsilon}^{p}\right)$ is relatively compact in the weak topology. Since $\gamma_{0}=\lim \gamma_{\varepsilon}$ is smooth on $\mathbb{C}^{n} \backslash S$ and depends only on $n-q$ variables $(n-q \leqslant p)$, we have $\lim \gamma_{\varepsilon}^{p}=\gamma_{0}^{p}=0$ on $\mathbb{C}^{n} \backslash S$. This shows that every weak limit of $\left(\gamma_{\varepsilon}^{p}\right)$ has support in $S$. Each of these is the direct image by inclusion of a unitary and homothety invariant $(p+q-n, p+q-n)$-form on $S$ with Lelong number equal to 1 at 0 . Therefore we must have

$$
\lim _{\varepsilon \rightarrow 0} \gamma_{\varepsilon}^{p}=\left(i_{S}\right)_{\star}\left(\alpha_{S}^{p+q-n}\right)=[S] \wedge \alpha_{S}^{p+q-n}
$$

and our claim is proved (of course, this can also be checked by direct elementary calculations). By taking the limsup in (7.14) we obtain

$$
\nu\left(T_{\lceil S}, 0, r+0\right)=\int_{\bar{B}(0, r)} T \wedge[S] \wedge \alpha_{S}^{p+q-n} \geqslant \nu(T, 0)
$$

(the singularity of $T$ at 0 does not create any difficulty because we can modify $T$ by a $d d^{c}$-exact form near 0 to make it smooth everywhere). Property b) follows when $r$ tends to 0 .

## § 8. Siu's Semicontinuity Theorem

Let $X, Y$ be complex manifolds of dimension $n, m$ such that $X$ is Stein. Let $\varphi$ : $X \times Y \longrightarrow[-\infty,+\infty[$ be a continuous plurisubharmonic function. We assume that $\varphi$ is semi-exhaustive with respect to Supp $T$, i.e. that for every compact subset $L \subset Y$ there exists $R=R(L)<0$ such that

$$
\begin{equation*}
\{(x, y) \in \operatorname{Supp} T \times L ; \varphi(x, y) \leqslant R\} \subset \subset X \times Y \tag{8.1}
\end{equation*}
$$

Let $T$ be a closed positive current of bidimension $(p, p)$ on $X$. For every point $y \in Y$, the function $\varphi_{y}(x):=\varphi(x, y)$ is semi-exhaustive on Supp $T$; one can therefore associate
with $y$ a generalized Lelong number $\nu\left(T, \varphi_{y}\right)$. Proposition 5.13 implies that the map $y \mapsto \nu\left(T, \varphi_{y}\right)$ is upper semi-continuous, hence the upperlevel sets

$$
\begin{equation*}
E_{c}=E_{c}(T, \varphi)=\left\{y \in Y ; \nu\left(T, \varphi_{y}\right) \geqslant c\right\}, c>0 \tag{8.2}
\end{equation*}
$$

are closed. Under mild additional hypotheses, we are going to show that the sets $E_{c}$ are in fact analytic subsets of $Y$, following [Demailly 1987a].
(8.3) Definition. We say that a function $f(x, y)$ is locally Hölder continuous with respect to $y$ on $X \times Y$ if every point of $X \times Y$ has a neighborhood $\Omega$ on which

$$
\left|f\left(x, y_{1}\right)-f\left(x, y_{2}\right)\right| \leqslant M\left|y_{1}-y_{2}\right|^{\gamma}
$$

for all $\left(x, y_{1}\right) \in \Omega,\left(x, y_{2}\right) \in \Omega$, with some constants $\left.\left.M>0, \gamma \in\right] 0,1\right]$, and suitable coordinates on $Y$.
(8.4) Theorem ([Demailly 1987a]). Let $T$ be a closed positive current on $X$ and let

$$
\varphi: X \times Y \longrightarrow[-\infty,+\infty[
$$

be a continuous plurisubharmonic function. Assume that $\varphi$ is semi-exhaustive on Supp $T$ and that $e^{\varphi(x, y)}$ is locally Hölder continuous with respect to $y$ on $X \times Y$. Then the upperlevel sets

$$
E_{c}(T, \varphi)=\left\{y \in Y ; \nu\left(T, \varphi_{y}\right) \geqslant c\right\}
$$

are analytic subsets of $Y$.

This theorem can be rephrased by saying that $y \longmapsto \nu\left(T, \varphi_{y}\right)$ is upper semi-continuous with respect to the analytic Zariski topology. As a special case, we get the following important result of [Siu 1974]:
(8.5) Corollary. If $T$ is a closed positive current of bidimension $(p, p)$ on a complex manifold $X$, the upperlevel sets $E_{c}(T)=\{x \in X ; \nu(T, x) \geqslant c\}$ of the usual Lelong numbers are analytic subsets of dimension $\leqslant p$.

Proof. The result is local, so we may assume that $X \subset \mathbb{C}^{n}$ is an open subset. Theorem 8.4 with $Y=X$ and $\varphi(x, y)=\log |x-y|$ shows that $E_{c}(T)$ is analytic. Moreover, Prop. 5.11 implies $\operatorname{dim} E_{c}(T) \leqslant p$.
(8.6) Generalization. Theorem 8.4 can be applied more generally to weight functions of the type

$$
\varphi(x, y)=\max _{j} \log \left(\sum_{k}\left|F_{j, k}(x, y)\right|^{\lambda_{j, k}}\right)
$$

where $F_{j, k}$ are holomorphic functions on $X \times Y$ and where $\gamma_{j, k}$ are positive real constants; in this case $e^{\varphi}$ is Hölder continuous of exponent $\gamma=\min \left\{\lambda_{j, k}, 1\right\}$ and $\varphi$ is semi-exhaustive with respect to the whole space $X$ as soon as the projection $\mathrm{pr}_{2}: \bigcap F_{j, k}^{-1}(0) \longrightarrow Y$ is proper and finite.

For example, when $\varphi(x, y)=\log \max \left|x_{j}-y_{j}\right|_{j}^{\lambda}$ on an open subset $X$ of $\mathbb{C}^{n}$, we see that the upperlevel sets for Kiselman's numbers $\nu(T, x, \lambda)$ are analytic in $X$ (a result first
proved in [Kiselman 1986]. More generally, set $\psi_{\lambda}(z)=\log \max \left|z_{j}\right|^{\lambda_{j}}$ and $\varphi(x, y, g)=$ $\psi_{\lambda}(g(x-y))$ where $x, y \in \mathbb{C}^{n}$ and $g \in \operatorname{Gl}\left(\mathbb{C}^{n}\right)$. Then $\nu\left(T, \varphi_{y, g}\right)$ is the Kiselman number of $T$ at $y$ when the coordinates have been rotated by $g$. It is clear that $\varphi$ is plurisubharmonic in $(x, y, g)$ and semi-exhaustive with respect to $x$, and that $e^{\varphi}$ is locally Hölder continuous with respect to $(y, g)$. Thus the upperlevel sets

$$
E_{c}=\left\{(y, g) \in X \times \operatorname{Gl}\left(\mathbb{C}^{n}\right) ; \nu\left(T, \varphi_{y, g}\right) \geqslant c\right\}
$$

are analytic in $X \times \mathrm{Gl}\left(\mathbb{C}^{n}\right)$. However this result is not meaningful on a manifold, because it is not invariant under coordinate changes. One can obtain an invariant version as follows. Let $X$ be a manifold and let $J^{k} \mathscr{O}_{X}$ be the bundle of $k$-jets of holomorphic functions on $X$. We consider the bundle $S_{k}$ over $X$ whose fiber $S_{k, y}$ is the set of $n$-tuples of $k$-jets $u=\left(u_{1}, \ldots, u_{n}\right) \in\left(J^{k} \mathcal{O}_{X, y}\right)^{n}$ such that $u_{j}(y)=0$ and $d u_{1} \wedge \ldots \wedge d u_{n}(y) \neq 0$. Let $\left(z_{j}\right)$ be local coordinates on an open set $\Omega \subset X$. Modulo $O\left(|z-y|^{k+1}\right)$, we can write

$$
u_{j}(z)=\sum_{1 \leqslant|\alpha| \leqslant k} a_{j, \alpha}(z-y)^{\alpha}
$$

with $\operatorname{det}\left(a_{j,\left(0, \ldots, 1_{k}, \ldots, 0\right)}\right) \neq 0$. The numbers $\left(\left(y_{j}\right),\left(a_{j, \alpha}\right)\right)$ define a coordinate system on the total space of $S_{k \upharpoonright \Omega}$. For $(x,(y, u)) \in X \times S_{k}$, we introduce the function

$$
\varphi(x, y, u)=\log \max \left|u_{j}(x)\right|^{\lambda_{j}}=\log \max \left|\sum_{1 \leqslant|\alpha| \leqslant k} a_{j, \alpha}(x-y)^{\alpha}\right|^{\lambda_{j}}
$$

which has all properties required by Th. 8.4 on a neighborhood of the diagonal $x=y$, i.e. a neighborhood of $X \times_{X} S_{k}$ in $X \times S_{k}$. For $k$ large, we claim that Kiselman's directional Lelong numbers

$$
\nu(T, y, u, \lambda):=\nu\left(T, \varphi_{y, u}\right)
$$

with respect to the coordinate system $\left(u_{j}\right)$ at $y$ do not depend on the selection of the jet representives and are therefore canonically defined on $S_{k}$. In fact, a change of $u_{j}$ by $O\left(|z-y|^{k+1}\right)$ adds $O\left(|z-y|^{(k+1) \lambda_{j}}\right)$ to $e^{\varphi}$, and we have $e^{\varphi} \geqslant O\left(|z-y|^{\max \lambda_{j}}\right)$. Hence by (7.1) it is enough to take $k+1 \geqslant \max \lambda_{j} / \min \lambda_{j}$. Theorem 8.4 then shows that the upperlevel sets $E_{c}(T, \varphi)$ are analytic in $S_{k}$.

Proof of the Semicontinuity Theorem 8.4. As the result is local on $Y$, we may assume without loss of generality that $Y$ is a ball in $\mathbb{C}^{m}$. After addition of a constant to $\varphi$, we may also assume that there exists a compact subset $K \subset X$ such that

$$
\{(x, y) \in X \times Y ; \varphi(x, y) \leqslant 0\} \subset K \times Y
$$

By Th. 7.1, the Lelong numbers depend only on the asymptotic behaviour of $\varphi$ near the (compact) polar set $\varphi^{-1}(-\infty) \cap(\operatorname{Supp} T \times Y)$. We can add a smooth strictly plurisubharmonic function on $X \times Y$ to make $\varphi$ strictly plurisuharmonic. Then Richberg's approximation theorem for continuous plurisubharmonic functions shows that there exists a smooth plurisubharmonic function $\widetilde{\varphi}$ such that $\varphi \leqslant \widetilde{\varphi} \leqslant \varphi+1$. We may therefore assume that $\varphi$ is smooth on $(X \times Y) \backslash \varphi^{-1}(-\infty)$.

- First step: construction of a local plurisubharmonic potential.

Our goal is to generalize the usual construction of plurisubharmonic potentials associated with a closed positive current ([Lelong 1967], [Skoda 1972a]). We replace here the usual kernel $|z-\zeta|^{-2 p}$ arising from the hermitian metric of $\mathbb{C}^{n}$ by a kernel depending on the weight $\varphi$. Let $\chi \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$ be an increasing function such that $\chi(t)=t$ for $t \leqslant-1$ and $\chi(t)=0$ for $t \geqslant 0$. We consider the half-plane $H=\{z \in \mathbb{C} ; \operatorname{Re} z<-1\}$ and associate with $T$ the potential function $V$ on $Y \times H$ defined by

$$
\begin{equation*}
V(y, z)=-\int_{\operatorname{Re} z}^{0} \nu\left(T, \varphi_{y}, t\right) \chi^{\prime}(t) d t . \tag{8.7}
\end{equation*}
$$

For every $t>\operatorname{Re} z$, Stokes' formula gives

$$
\nu\left(T, \varphi_{y}, t\right)=\int_{\varphi(x, y)<t} T(x) \wedge\left(d d_{x}^{c} \widetilde{\varphi}(x, y, z)\right)^{p}
$$

with $\widetilde{\varphi}(x, y, z):=\max \{\varphi(x, y), \operatorname{Re} z\}$. The Fubini theorem applied to (8.7) gives

$$
\begin{aligned}
V(y, z) & =-\int_{\substack{x \in X, \varphi(x, y)<t \\
\operatorname{Re} z<t<0}} T(x) \wedge\left(d d_{x}^{c} \widetilde{\varphi}(x, y, z)\right)^{p} \chi^{\prime}(t) d t \\
& =\int_{x \in X} T(x) \wedge \chi(\widetilde{\varphi}(x, y, z))\left(d d_{x}^{c} \widetilde{\varphi}(x, y, z)\right)^{p} .
\end{aligned}
$$

For all ( $n-1, n-1$ )-form $h$ of class $C^{\infty}$ with compact support in $Y \times H$, we get

$$
\begin{aligned}
\left\langle d d^{c} V, h\right\rangle & =\left\langle V, d d^{c} h\right\rangle \\
& =\int_{X \times Y \times H} T(x) \wedge \chi(\widetilde{\varphi}(x, y, z))\left(d d^{c} \widetilde{\varphi}(x, y, z)\right)^{p} \wedge d d^{c} h(y, z) .
\end{aligned}
$$

Observe that the replacement of $d d_{x}^{c}$ by the total differentiation $d d^{c}=d d_{x, y, z}^{c}$ does not modify the integrand, because the terms in $d x, d \bar{x}$ must have total bidegree $(n, n)$. The current $T(x) \wedge \chi(\widetilde{\varphi}(x, y, z)) h(y, z)$ has compact support in $X \times Y \times H$. An integration by parts can thus be performed to obtain

$$
\left\langle d d^{c} V, h\right\rangle=\int_{X \times Y \times H} T(x) \wedge d d^{c}(\chi \circ \widetilde{\varphi}(x, y, z)) \wedge\left(d d^{c} \widetilde{\varphi}(x, y, z)\right)^{p} \wedge h(y, z) .
$$

On the corona $\{-1 \leqslant \varphi(x, y) \leqslant 0\}$ we have $\widetilde{\varphi}(x, y, z)=\varphi(x, y)$, whereas for $\varphi(x, y)<-1$ we get $\widetilde{\varphi}<-1$ and $\chi \circ \widetilde{\varphi}=\widetilde{\varphi}$. As $\widetilde{\varphi}$ is plurisubharmonic, we see that $d d^{c} V(y, z)$ is the sum of the positive (1,1)-form

$$
(y, z) \longmapsto \int_{\{x \in X ; \varphi(x, y)<-1\}} T(x) \wedge\left(d d_{x, y, z}^{c} \widetilde{\varphi}(x, y, z)\right)^{p+1}
$$

and of the $(1,1)$-form independent of $z$

$$
y \longmapsto \int_{\{x \in X ;-1 \leqslant \varphi(x, y) \leqslant 0\}} T \wedge d d_{x, y}^{c}(\chi \circ \varphi) \wedge\left(d d_{x, y}^{c} \varphi\right)^{p} .
$$

As $\varphi$ is smooth outside $\varphi^{-1}(-\infty)$, this last form has locally bounded coefficients. Hence $d d^{c} V(y, z)$ is $\geqslant 0$ except perhaps for locally bounded terms. In addition, $V$ is continuous
on $Y \times H$ because $T \wedge\left(d d^{c} \widetilde{\varphi}_{y, z}\right)^{p}$ is weakly continuous in the variables $(y, z)$ by Th. 3.5. We therefore obtain the following result.
(8.8) Proposition. There exists a positive plurisubharmonic function $\rho$ in $\mathscr{C}^{\infty}(Y)$ such that $\rho(y)+V(y, z)$ is plurisubharmonic on $Y \times H$.

If we let $\operatorname{Re} z$ tend to $-\infty$, we see that the function

$$
U_{0}(y)=\rho(y)+V(y,-\infty)=\rho(y)-\int_{-\infty}^{0} \nu\left(T, \varphi_{y}, t\right) \chi^{\prime}(t) d t
$$

is locally plurisubharmonic or $\equiv-\infty$ on $Y$. Furthermore, it is clear that $U_{0}(y)=-\infty$ at every point $y$ such that $\nu\left(T, \varphi_{y}\right)>0$. If $Y$ is connected and $U_{0} \not \equiv-\infty$, we already conclude that the density set $\bigcup_{c>0} E_{c}$ is pluripolar in $Y$.

- Second step: application of Kiselman's minimum principle.

Let $a \geqslant 0$ be arbitrary. The function

$$
Y \times H \ni(y, z) \longmapsto \rho(y)+V(y, z)-a \operatorname{Re} z
$$

is plurisubharmonic and independent of $\operatorname{Im} z$. By Kiselman's theorem 1.7.8, the Legendre transform

$$
U_{a}(y)=\inf _{r<-1}\{\rho(y)+V(y, r)-a r\}
$$

is locally plurisubharmonic or $\equiv-\infty$ on $Y$.
(8.9) Lemma. Let $y_{0} \in Y$ be a given point.
a) If $a>\nu\left(T, \varphi_{y_{0}}\right)$, then $U_{a}$ is bounded below on a neighborhood of $y_{0}$.
b) If $a<\nu\left(T, \varphi_{y_{0}}\right)$, then $U_{a}\left(y_{0}\right)=-\infty$.

Proof. By definition of $V$ (cf. (8.7)) we have

$$
\begin{equation*}
V(y, r) \leqslant-\nu\left(T, \varphi_{y}, r\right) \int_{r}^{0} \chi^{\prime}(t) d t=r \nu\left(T, \varphi_{y}, r\right) \leqslant r \nu\left(T, \varphi_{y}\right) \tag{8.10}
\end{equation*}
$$

Then clearly $U_{a}\left(y_{0}\right)=-\infty$ if $a<\nu\left(T, \varphi_{y_{0}}\right)$. On the other hand, if $\nu\left(T, \varphi_{y_{0}}\right)<a$, there exists $t_{0}<0$ such that $\nu\left(T, \varphi_{y_{0}}, t_{0}\right)<a$. Fix $r_{0}<t_{0}$. The semi-continuity property (5.13) shows that there exists a neighborhood $\omega$ of $y_{0}$ such that $\sup _{y \in \omega} \nu\left(T, \varphi_{y}, r_{0}\right)<a$. For all $y \in \omega$, we get

$$
V(y, r) \geqslant-C-a \int_{r}^{r_{0}} \chi^{\prime}(t) d t=-C+a\left(r-r_{0}\right)
$$

and this implies $U_{a}(y) \geqslant-C-a r_{0}$.
(8.11) Theorem. If $Y$ is connected and if $E_{c} \neq Y$, then $E_{c}$ is a closed complete pluripolar subset of $Y$, i.e. there exists a continuous plurisubharmonic function $w: Y \longrightarrow$ $\left[-\infty,+\infty\left[\right.\right.$ such that $E_{c}=w^{-1}(-\infty)$.

Proof. We first observe that the family $\left(U_{a}\right)$ is increasing in $a$, that $U_{a}=-\infty$ on $E_{c}$ for all $a<c$ and that $\sup _{a<c} U_{a}(y)>-\infty$ if $y \in Y \backslash E_{c}$ (apply Lemma 8.9). For any integer $k \geqslant 1$, let $w_{k} \in \mathscr{C}^{\infty}(Y)$ be a plurisubharmonic regularization of $U_{c-1 / k}$ such that $w_{k} \geqslant U_{c-1 / k}$ on $Y$ and $w_{k} \leqslant-2^{k}$ on $E_{c} \cap Y_{k}$ where $Y_{k}=\{y \in Y ; d(y, \partial Y) \geqslant 1 / k\}$. Then Lemma 8.9 a) shows that the family $\left(w_{k}\right)$ is uniformly bounded below on every compact subset of $Y \backslash E_{c}$. We can also choose $w_{k}$ uniformly bounded above on every compact subset of $Y$ because $U_{c-1 / k} \leqslant U_{c}$. The function

$$
w=\sum_{k=1}^{+\infty} 2^{-k} w_{k}
$$

satifies our requirements.

- Third step: estimation of the singularities of the potentials $U_{a}$.
(8.12) Lemma. Let $y_{0} \in Y$ be a given point, $L$ a compact neighborhood of $y_{0}, K \subset X$ a compact subset and $r_{0}$ a real number $<-1$ such that

$$
\left\{(x, y) \in X \times L ; \varphi(x, y) \leqslant r_{0}\right\} \subset K \times L
$$

Assume that $e^{\varphi(x, y)}$ is locally Hölder continuous in $y$ and that

$$
\left|f\left(x, y_{1}\right)-f\left(x, y_{2}\right)\right| \leqslant M\left|y_{1}-y_{2}\right|^{\gamma}
$$

for all $\left(x, y_{1}, y_{2}\right) \in K \times L \times L$. Then, for all $\left.\varepsilon \in\right] 0,1[$, there exists a real number $\eta(\varepsilon)>0$ such that all $y \in Y$ with $\left|y-y_{0}\right|<\eta(\varepsilon)$ satisfy

$$
U_{a}(y) \leqslant \rho(y)+\left((1-\varepsilon)^{p} \nu\left(T, \varphi_{y_{0}}\right)-a\right)\left(\gamma \log \left|y-y_{0}\right|+\log \frac{2 e M}{\varepsilon}\right) .
$$

Proof. First, we try to estimate $\nu\left(T, \varphi_{y}, r\right)$ when $y \in L$ is near $y_{0}$. Set

$$
\left\{\begin{array}{llr}
\psi(x)=(1-\varepsilon) \varphi_{y_{0}}(x)+\varepsilon r-\varepsilon / 2 & \text { if } & \varphi_{y_{0}}(x) \leqslant r-1 \\
\psi(x)=\max \left(\varphi_{y}(x),(1-\varepsilon) \varphi_{y_{0}}(x)+\varepsilon r-\varepsilon / 2\right) & \text { if } r-1 \leqslant \varphi_{y_{0}}(x) \leqslant r \\
\psi(x)=\varphi_{y}(x) & \text { if } & r \leqslant \varphi_{y_{0}}(x) \leqslant r_{0}
\end{array}\right.
$$

and verify that this definition is coherent when $\left|y-y_{0}\right|$ is small enough. By hypothesis

$$
\left|e^{\varphi_{y}(x)}-e^{\varphi_{y_{0}}(x)}\right| \leqslant M\left|y-y_{0}\right|^{\gamma} .
$$

This inequality implies

$$
\begin{aligned}
& \varphi_{y}(x) \leqslant \varphi_{y_{0}}(x)+\log \left(1+M\left|y-y_{0}\right|^{\gamma} e^{-\varphi_{y_{0}}(x)}\right) \\
& \varphi_{y}(x) \geqslant \varphi_{y_{0}}(x)+\log \left(1-M\left|y-y_{0}\right|^{\gamma} e^{-\varphi_{y_{0}}(x)}\right)
\end{aligned}
$$

In particular, for $\varphi_{y_{0}}(x)=r$, we have $(1-\varepsilon) \varphi_{y_{0}}(x)+\varepsilon r-\varepsilon / 2=r-\varepsilon / 2$, thus

$$
\varphi_{y}(x) \geqslant r+\log \left(1-M\left|y-y_{0}\right|^{\gamma} e^{-r}\right) .
$$

Similarly, for $\varphi_{y_{0}}(x)=r-1$, we have $(1-\varepsilon) \varphi_{y_{0}}(x)+\varepsilon r-\varepsilon / 2=r-1+\varepsilon / 2$, thus

$$
\varphi_{y}(x) \leqslant r-1+\log \left(1+M\left|y-y_{0}\right|^{\gamma} e^{1-r}\right) .
$$

The definition of $\psi$ is thus coherent as soon as $M\left|y-y_{0}\right|^{\gamma} e^{1-r} \leqslant \varepsilon / 2$, i.e.

$$
\gamma \log \left|y-y_{0}\right|+\log \frac{2 e M}{\varepsilon} \leqslant r .
$$

In this case $\psi$ coincides with $\varphi_{y}$ on a neighborhood of $\{\psi=r\}$, and with

$$
(1-\varepsilon) \varphi_{y_{0}}(x)+\varepsilon r-\varepsilon / 2
$$

on a neighborhood of the polar set $\psi^{-1}(-\infty)$. By Stokes' formula applied to $\nu(T, \psi, r)$, we infer

$$
\nu\left(T, \varphi_{y}, r\right)=\nu(T, \psi, r) \geqslant \nu(T, \psi)=(1-\varepsilon)^{p} \nu\left(T, \varphi_{y_{0}}\right) .
$$

From (8.10) we get $V(y, r) \leqslant r \nu\left(T, \varphi_{y}, r\right)$, hence

$$
\begin{align*}
& U_{a}(y) \leqslant \rho(y)+V(y, r)-a r \leqslant \rho(y)+r\left(\nu\left(T, \varphi_{y}, r\right)-a\right), \\
& U_{a}(y) \leqslant \rho(y)+r\left((1-\varepsilon)^{p} \nu\left(T, \varphi_{y_{0}}\right)-a\right) . \tag{8.13}
\end{align*}
$$

Suppose $\gamma \log \left|y-y_{0}\right|+\log (2 e M / \varepsilon) \leqslant r_{0}$, i.e. $\left|y-y_{0}\right| \leqslant(\varepsilon / 2 e M)^{1 / \gamma} e^{r_{0} / \gamma}$; one can then choose $r=\gamma \log \left|y-y_{0}\right|+\log (2 e M / \varepsilon)$, and by (8.13) this yields the inequality asserted in Th. 8.12.

- Fourth step: application of the Hörmander-Bombieri-Skoda theorem.

The end of the proof relies on the following crucial result, which is a consequence of the Hörmander-Bombieri-Skoda theorem ([Bombieri 1970], [Skoda 1972a], [Skoda 1976]); it will be proved in Chapter 8, see Cor. 8.?.?.
(8.14) Proposition. Let $u$ be a plurisubharmonic function on a complex manifold $Y$. The set of points in a neighborhood of which $e^{-u}$ is not integrable is an analytic subset of $Y$.

Proof of Theorem 8.4 (end).. The main idea in what follows is due to [Kiselman 1979]. For $a, b>0$, we let $Z_{a, b}$ be the set of points in a neighborhood of which $\exp \left(-U_{a} / b\right)$ is not integrable. Then $Z_{a, b}$ is analytic, and as the family $\left(U_{a}\right)$ is increasing in $a$, we have $Z_{a^{\prime}, b^{\prime}} \supset Z_{a^{\prime \prime}, b^{\prime \prime}}$ if $a^{\prime} \leqslant a^{\prime \prime}, b^{\prime} \leqslant b^{\prime \prime}$.

Let $y_{0} \in Y$ be a given point. If $y_{0} \notin E_{c}$, then $\nu\left(T, \varphi_{y_{0}}\right)<c$ by definition of $E_{c}$. Choose $a$ such that $\nu\left(T, \varphi_{y_{0}}\right)<a<c$. Lemma 8.9 a) implies that $U_{a}$ is bounded below in a neighborhood of $y_{0}$, thus $\exp \left(-U_{a} / b\right)$ is integrable and $y_{0} \notin Z_{a, b}$ for all $b>0$.

On the other hand, if $y_{0} \in E_{c}$ and if $a<c$, then Lemma 8.12 implies for all $\varepsilon>0$ that

$$
U_{a}(y) \leqslant(1-\varepsilon)(c-a) \gamma \log \left|y-y_{0}\right|+C(\varepsilon)
$$

on a neighborhood of $y_{0}$. Hence $\exp \left(-U_{a} / b\right)$ is non integrable at $y_{0}$ as soon as $b<$ $(c-a) \gamma / 2 m$, where $m=\operatorname{dim} Y$. We obtain therefore

$$
E_{c}=\bigcap_{\substack{a<c \\ b<(c-a) \gamma / 2 m}} Z_{a, b} .
$$

This proves that $E_{c}$ is an analytic subset of $Y$.
Finally, we use Cor. 8.5 to derive an important decomposition formula for currents, which is again due to [Siu 1974]. We first begin by two simple observations.
(8.15) Lemma. If $T$ is a closed positive current of bidimension ( $p, p$ ) and $A$ is an irreducible analytic set in $X$, we set

$$
m_{A}=\inf \{\nu(T, x) ; x \in A\}
$$

Then $\nu(T, x)=m_{A}$ for all $x \in A \backslash \bigcup A_{j}^{\prime}$, where $\left(A_{j}^{\prime}\right)$ is a countable family of proper analytic subsets of $A$. We say that $m_{A}$ is the generic Lelong number of $T$ along $A$.

Proof. By definition of $m_{A}$ and $E_{c}(T)$, we have $\nu(T, x) \geqslant m_{A}$ for every $x \in A$ and

$$
\nu(T, x)=m_{A} \quad \text { on } \quad A \backslash \bigcup_{c \in \mathbb{Q}, c>m_{A}} A \cap E_{c}(T) .
$$

However, for $c>m_{A}$, the intersection $A \cap E_{c}(T)$ is a proper analytic subset of $A$.
(8.16) Proposition. Let $T$ be a closed positive current of bidimension ( $p, p$ ) and let $A$ be an irreducible $p$-dimensional analytic subset of $X$. Then $\mathbb{1}_{A} T=m_{A}[A]$, in particular $T-m_{A}[A]$ is positive.

Proof. As the question is local and as a closed positive current of bidimension ( $p, p$ ) cannot carry any mass on a ( $p-1$ )-dimensional analytic subset, it is enough to work in a neighborhood of a regular point $x_{0} \in A$. Hence, by choosing suitable coordinates, we can suppose that $X$ is an open set in $\mathbb{C}^{n}$ and that $A$ is the intersection of $X$ with a $p$-dimensional linear subspace. Then, for every point $a \in A$, the inequality $\nu(T, a) \geqslant m_{A}$ implies

$$
\sigma_{T}(B(a, r)) \geqslant m_{A} \pi^{p} r^{2 p} / p!=m_{A} \sigma_{[A]}(B(a, r))
$$

for all $r$ such that $B(a, r) \subset X$. Now, set $\Theta=T-m_{A}[A]$ and $\beta=d d^{c}|z|^{2}$. Our inequality says that $\int \mathbb{1}_{B(a, r)} \Theta \wedge \beta^{p} \geqslant 0$. If we integrate this with respect to some positive continuous function $f$ with compact support in $A$, we get $\int_{X} g_{r} \Theta \wedge \beta^{p} \geqslant 0$ where

$$
g_{r}(z)=\int_{A} \mathbb{1}_{B(a, r)}(z) f(a) d \lambda(a)=\int_{a \in A \cap B(z, r)} f(a) d \lambda(a) .
$$

Here $g_{r}$ is continuous on $\mathbb{C}^{n}$, and as $r$ tends to 0 the function $g_{r}(z) /\left(\pi^{p} r^{2 p} / p!\right)$ converges to $f$ on $A$ and to 0 on $X \backslash A$, with a global uniform bound. Hence we obtain $\int \mathbb{1}_{A} f \Theta \wedge \beta^{p} \geqslant$ 0 . Since this inequality is true for all continuous functions $f \geqslant 0$ with compact support in $A$, we conclude that the measure $\mathbb{1}_{A} \Theta \wedge \beta^{p}$ is positive. By a linear change of coordinates, we see that

$$
\mathbb{1}_{A} \Theta \wedge\left(d d^{c} \sum_{1 \leqslant j \leqslant n} \lambda_{j}\left|u_{j}\right|^{2}\right)^{n} \geqslant 0
$$

for every basis $\left(u_{1}, \ldots, u_{n}\right)$ of linear forms and for all coefficients $\lambda_{j}>0$. Take $\lambda_{1}=\ldots=$ $\lambda_{p}=1$ and let the other $\lambda_{j}$ tend to 0 . Then we get

$$
\mathbb{1}_{A} \Theta \wedge \mathrm{i} d u_{1} \wedge d \bar{u}_{1} \wedge \ldots \wedge d u_{p} \wedge d \bar{u}_{p} \geqslant 0
$$

This implies $\mathbb{1}_{A} \Theta \geqslant 0$, or equivalently $\mathbb{1}_{A} T \geqslant m_{A}[A]$. By Cor. 2.4 we know that $\mathbb{1}_{A} T$ is a closed positive current, thus $\mathbb{1}_{A} T=\lambda[A]$ with $\lambda \geqslant 0$. We have just seen that $\lambda \geqslant m_{A}$. On the other hand, $T \geqslant \mathbb{1}_{A} T=\lambda[A]$ clearly implies $m_{A} \geqslant \lambda$.
(8.16) Siu's decomposition formula. If $T$ is a closed positive current of bidimension ( $p, p$ ), there is a unique decomposition of $T$ as a (possibly finite) weakly convergent series

$$
T=\sum_{j \geqslant 1} \lambda_{j}\left[A_{j}\right]+R, \quad \lambda_{j}>0,
$$

where $\left[A_{j}\right]$ is the current of integration over an irreducible $p$-dimensional analytic set $A_{j} \subset X$ and where $R$ is a closed positive current with the property that $\operatorname{dim} E_{c}(R)<p$ for every $c>0$.

Proof of uniqueness.. If $T$ has such a decomposition, the $p$-dimensional components of $E_{c}(T)$ are $\left(A_{j}\right)_{\lambda_{j} \geqslant c}$, for $\nu(T, x)=\sum \lambda_{j} \nu\left(\left[A_{j}\right], x\right)+\nu(R, x)$ is non zero only on $\bigcup A_{j} \cup$ $\bigcup E_{c}(R)$, and is equal to $\lambda_{j}$ generically on $A_{j}$ (more precisely, $\nu(T, x)=\lambda_{j}$ at every regular point of $A_{j}$ which does not belong to any intersection $A_{j} \cup A_{k}, k \neq j$ or to $\bigcup E_{c}(R)$ ). In particular $A_{j}$ and $\lambda_{j}$ are unique.
Proof of existence.. Let $\left(A_{j}\right)_{j \geqslant 1}$ be the countable collection of $p$-dimensional components occurring in one of the sets $E_{c}(T), c \in \mathbb{Q}_{+}^{\star}$, and let $\lambda_{j}>0$ be the generic Lelong number of $T$ along $A_{j}$. Then Prop. 8.16 shows by induction on $N$ that $R_{N}=T-\sum_{1 \leqslant j \leqslant N} \lambda_{j}\left[A_{j}\right]$ is positive. As $R_{N}$ is a decreasing sequence, there must be a limit $R=\lim _{N \rightarrow+\infty} R_{N}$ in the weak topology. Thus we have the asserted decomposition. By construction, $R$ has zero generic Lelong number along $A_{j}$, so $\operatorname{dim} E_{c}(R)<p$ for every $c>0$.

It is very important to note that some components of lower dimension can actually occur in $E_{c}(R)$, but they cannot be subtracted because $R$ has bidimension $(p, p)$. A typical case is the case of a bidimension $(n-1, n-1)$ current $T=d d^{c} u$ with $u=$ $\log \left(\left|F_{j}\right|^{\gamma_{1}}+\ldots\left|F_{N}\right|^{\gamma_{N}}\right)$ and $F_{j} \in \mathscr{O}(X)$. In general $\bigcup E_{c}(T)=\bigcap F_{j}^{-1}(0)$ has dimension $<$ $n-1$. In that case, an important formula due to King plays the role of (8.17). We state it in a somewhat more general form than its original version [King 1970].
(8.18) King's formula. Let $F_{1}, \ldots, F_{N}$ be holomorphic functions on a complex manifold $X$, such that the zero variety $Z=\bigcap F_{j}^{-1}(0)$ has codimension $\geqslant p$, and set $u=$ $\log \sum\left|F_{j}\right|^{\gamma_{j}}$ with arbitrary coefficients $\gamma_{j}>0$. Let $\left(Z_{k}\right)_{k \geqslant 1}$ be the irreducible components of $Z$ of codimension $p$ exactly. Then there exist multiplicities $\lambda_{k}>0$ such that

$$
\left(d d^{c} u\right)^{p}=\sum_{k \geqslant 1} \lambda_{k}\left[Z_{k}\right]+R,
$$

where $R$ is a closed positive current such that $\mathbb{1}_{Z} R=0$ and $\operatorname{codim} E_{c}(R)>p$ for every $c>0$. Moreover the multiplicities $\lambda_{k}$ are integers if $\gamma_{1}, \ldots, \gamma_{N}$ are integers, and $\lambda_{k}=$ $\gamma_{1} \ldots \gamma_{p}$ if $\gamma_{1} \leqslant \ldots \leqslant \gamma_{N}$ and some partial Jacobian determinant of $\left(F_{1}, \ldots, F_{p}\right)$ of order $p$ does not vanish identically along $Z_{k}$.

Proof. Observe that $\left(d d^{c} u\right)^{p}$ is well defined thanks to Cor. 4.11. The comparison theorem 7.8 applied with $\varphi(z)=\log |z-x|, v_{1}=\ldots=v_{p}=u, u_{1}=\ldots=u_{p}=\varphi$ and $T=1$ shows that the Lelong number of $\left(d d^{c} u\right)^{p}$ is equal to 0 at every point of $X \backslash Z$. Hence $E_{c}\left(\left(d d^{c} u\right)^{p}\right)$ is contained in $Z$ and its $(n-p)$-dimensional components are members of the
family $\left(Z_{k}\right)$. The asserted decomposition follows from Siu's formula 8.16. We must have $\mathbb{1}_{Z_{k}} R=0$ for all irreducible components of $Z$ : when $\operatorname{codim} Z_{k}>p$ this is automatically true, and when $\operatorname{codim} Z_{k}=p$ this follows from (8.16) and the fact that $\operatorname{codim} E_{c}(R)>p$. If $\operatorname{det}\left(\partial F_{j} / \partial z_{k}\right)_{1 \leqslant j, k \leqslant p} \neq 0$ at some point $x_{0} \in Z_{k}$, then $\left(Z, x_{0}\right)=\left(Z_{k}, x_{0}\right)$ is a smooth germ defined by the equations $F_{1}=\ldots=F_{p}=0$. If we denote $v=\log \sum_{j \leqslant p}\left|F_{j}\right|^{\gamma_{j}}$ with $\gamma_{1} \leqslant \ldots \leqslant \gamma_{N}$, then $u \sim v$ near $Z_{k}$ and Th. 7.8 implies $\nu\left(\left(d d^{c} u\right)^{p}, x\right)=\nu\left(\left(d d^{c} v\right)^{p}, x\right)$ for all $x \in Z_{k}$ near $x_{0}$. On the other hand, if $G:=\left(F_{1}, \ldots, F_{p}\right): X \rightarrow \mathbb{C}^{p}$, Cor. 7.4 gives

$$
\left(d d^{c} v\right)^{p}=G^{\star}\left(d d^{c} \log \sum_{1 \leqslant j \leqslant p}\left|z_{j}\right|^{\gamma_{j}}\right)^{p}=\gamma_{1} \ldots \gamma_{p} G^{\star} \delta_{0}=\gamma_{1} \ldots \gamma_{p}\left[Z_{k}\right]
$$

near $x_{0}$. This implies that the generic Lelong number of $\left(d d^{c} u\right)^{p}$ along $Z_{k}$ is $\lambda_{k}=\gamma_{1} \ldots \gamma_{p}$. The integrality of $\lambda_{k}$ when $\gamma_{1}, \ldots, \gamma_{N}$ are integers will be proved in the next section.

## § 9. Transformation of Lelong Numbers by Direct Images

Let $F: X \rightarrow Y$ be a holomorphic map between complex manifolds of respective dimensions $\operatorname{dim} X=n, \operatorname{dim} Y=m$, and let $T$ be a closed positive current of bidimension $(p, p)$ on $X$. If $F_{\uparrow \text { Supp } T}$ is proper, the direct image $F_{\star} T$ is defined by

$$
\begin{equation*}
\left\langle F_{\star} T, \alpha\right\rangle=\left\langle T, F^{\star} \alpha\right\rangle \tag{9.1}
\end{equation*}
$$

for every test form $\alpha$ of bidegree $(p, p)$ on $Y$. This makes sense because Supp $T \cap$ $F^{-1}(\operatorname{Supp} \alpha)$ is compact. It is easily seen that $F_{\star} T$ is a closed positive current of bidimension ( $p, p$ ) on $Y$.
(9.2) Example. Let $T=[A]$ where $A$ is a $p$-dimensional irreducible analytic set in $X$ such that $F_{\lceil A}$ is proper. We know by Remmert's theorem 2.7.8 that $F(A)$ is an analytic set in $Y$. Two cases may occur. Either $F_{\lceil A}$ is generically finite and $F$ induces an étale covering $A \backslash F^{-1}(Z) \longrightarrow F(A) \backslash Z$ for some nowhere dense analytic subset $Z \subset F(A)$, or $F_{\upharpoonright A}$ has generic fibers of positive dimension and $\operatorname{dim} F(A)<\operatorname{dim} A$. In the first case, let $s<+\infty$ be the covering degree. Then for every test form $\alpha$ of bidegree $(p, p)$ on $Y$ we get

$$
\left\langle F_{\star}[A], \alpha\right\rangle=\int_{A} F^{\star} \alpha=\int_{A \backslash F^{-1}(Z)} F^{\star} \alpha=s \int_{F(A) \backslash Z} \alpha=s\langle[F(A)], \alpha\rangle
$$

because $Z$ and $F^{-1}(Z)$ are negligible sets. Hence $F_{\star}[A]=s[F(A)]$. On the other hand, if $\operatorname{dim} F(A)<\operatorname{dim} A=p$, the restriction of $\alpha$ to $F(A)_{\text {reg }}$ is zero, and therefore so is this the restriction of $F^{\star} \alpha$ to $A_{\text {reg }}$. Hence $F_{\star}[A]=0$.

Now, let $\psi$ be a continuous plurisubharmonic function on $Y$ which is semi-exhaustive on $F(\operatorname{Supp} T)$ (this set certainly contains $\left.\operatorname{Supp} F_{\star} T\right)$. Since $F_{\upharpoonright \operatorname{Supp} T}$ is proper, it follows that $\psi \circ F$ is semi-exhaustive on Supp $T$, for

$$
\text { Supp } T \cap\{\psi \circ F<R\}=F^{-1}(F(\operatorname{Supp} T) \cap\{\psi<R\})
$$

(9.3) Proposition. If $F(\operatorname{Supp} T) \cap\{\psi<R\} \subset \subset Y$, we have

$$
\nu\left(F_{\star} T, \psi, r\right)=\nu(T, \psi \circ F, r) \quad \text { for all } r<R,
$$

in particular $\nu\left(F_{\star} T, \psi\right)=\nu(T, \psi \circ F)$.

Here, we do not necessarily assume that $X$ or $Y$ are Stein; we thus replace $\psi$ with $\psi_{\geqslant s}=\max \{\psi, s\}, s<r$, in the definition of $\nu\left(F_{\star} T, \psi, r\right)$ and $\nu(T, \psi \circ F, r)$.

Proof. The first equality can be written

$$
\int_{Y} F_{\star} T \wedge \mathbb{1}_{\{\psi<r\}}\left(d d^{c} \psi_{\geqslant s}\right)^{p}=\int_{X} T \wedge\left(\mathbb{1}_{\{\psi<r\}} \circ F\right)\left(d d^{c} \psi_{\geqslant s} \circ F\right)^{p}
$$

This follows almost immediately from the adjunction formula (9.1) when $\psi$ is smooth and when we write $\mathbb{1}_{\{\psi<R\}}=\lim \uparrow g_{k}$ for some sequence of smooth functions $g_{k}$. In general, we write $\psi_{\geqslant s}$ as a decreasing limit of smooth plurisubharmonic functions and we apply our monotone continuity theorems (if $Y$ is not Stein, Richberg's theorem shows that we can obtain a decreasing sequence of almost plurisubharmonic approximations such that the negative part of $d d^{c}$ converges uniformly to 0 ; this is good enough to apply the monotone continuity theorem; note that the integration is made on compact subsets, thanks to the semi-exhaustivity assumption on $\psi$ ).

It follows from this that understanding the transformation of Lelong numbers under direct images is equivalent to understanding the effect of $F$ on the weight. We are mostly interested in computing the ordinary Lelong numbers $\nu\left(F_{\star} T, y\right)$ associated with the weight $\psi(w)=\log |w-y|$ in some local coordinates $\left(w_{1}, \ldots, w_{m}\right)$ on $Y$ near $y$. Then Prop. 9.3 gives

$$
\begin{align*}
\nu\left(F_{\star} T, y\right) & =\nu(T, \log |F-y|) \quad \text { with }  \tag{9.4}\\
\log |F(z)-y| & =\frac{1}{2} \log \sum\left|F_{j}(z)-y_{j}\right|^{2}, \quad F_{j}=w_{j} \circ F .
\end{align*}
$$

We are going to show that $\nu(T, \log |F-y|)$ is bounded below by a linear combination of the Lelong numbers of $T$ at points $x$ in the fiber $F^{-1}(y)$, with suitable multiplicities attached to $F$ at these points. These multiplicities can be seen as generalizations of the notion of multiplicity of an analytic map introduced by [Stoll 1966].
(9.5) Definition. Let $x \in X$ and $y=F(x)$. Suppose that the codimension of the fiber $F^{-1}(y)$ at $x$ is $\geqslant p$. Then we set

$$
\mu_{p}(F, x)=\nu\left(\left(d d^{c} \log |F-y|\right)^{p}, x\right) .
$$

Observe that $\left(d d^{c} \log |F-y|\right)^{p}$ is well defined thanks to Cor. 4.10. The second comparison theorem 7.8 immediately shows that $\mu_{p}(F, x)$ is independent of the choice of local coordinates on $Y$ (and also on $X$, since Lelong nombers do not depend on coordinates). By definition, $\mu_{p}(F, x)$ is the mass carried by $\{x\}$ of the measure

$$
\left(d d^{c} \log |F(z)-y|\right)^{p} \wedge\left(d d^{c} \log |z-x|\right)^{n-p} .
$$

We are going to give a more geometric interpretation of this multiplicity, from which it will follow that $\mu_{p}(F, x)$ is always a positive integer (in particular, the proof of (8.18) will be complete).
(9.6) Example. For $p=n=\operatorname{dim} X$, the assumption $\operatorname{codim}_{x} F^{-1}(y) \geqslant p$ means that the germ of map $F:(X, x) \longrightarrow(Y, y)$ is finite. Let $U_{x}$ be a neighborhood of $x$ such that $\bar{U}_{x} \cap F^{-1}(y)=\{x\}$, let $W_{y}$ be a neighborhood of $y$ disjoint from $F\left(\partial U_{x}\right)$ and let $V_{x}=U_{x} \cap F^{-1}\left(W_{y}\right)$. Then $F: V_{x} \rightarrow W_{y}$ is proper and finite, and we have $F_{\star}\left[V_{x}\right]=s\left[F\left(V_{x}\right)\right]$ where $s$ is the local covering degree of $F: V_{x} \rightarrow F\left(V_{x}\right)$ at $x$. Therefore

$$
\begin{aligned}
\mu_{n}(F, x) & =\int_{\{x\}}\left(d d^{c} \log |F-y|\right)^{n}=\nu\left(\left[V_{x}\right], \log |F-y|\right)=\nu\left(F_{\star}\left[V_{x}\right], y\right) \\
& =s \nu\left(F\left(V_{x}\right), y\right)
\end{aligned}
$$

In the particular case when $\operatorname{dim} Y=\operatorname{dim} X$, we have $\left(F\left(V_{x}\right), y\right)=(Y, y)$, so $\mu_{n}(F, x)=s$. In general, it is a well known fact that the ideal generated by $\left(F_{1}-y_{1}, \ldots, F_{m}-y_{m}\right)$ in $0_{X, x}$ has the same integral closure as the ideal generated by $n$ generic linear combinations of the generators, that is, for a generic choice of coordinates $w^{\prime}=\left(w_{1}, \ldots, w_{n}\right), w^{\prime \prime}=$ $\left(w_{n+1}, \ldots, w_{m}\right)$ on $(Y, y)$, we have $|F(z)-y| \leqslant C\left|w^{\prime} \circ F(z)\right|$ (this is a simple consequence of Lemma 7.5 applied to $A=F\left(V_{x}\right)$ ). Hence for $p=n$, the comparison theorem 7.1 gives

$$
\mu_{n}(F, x)=\mu_{n}\left(w^{\prime} \circ F, x\right)=\text { local covering degree of } w^{\prime} \circ F \text { at } x,
$$

for a generic choice of coordinates $\left(w^{\prime}, w^{\prime \prime}\right)$ on $(Y, y)$.
(9.7) Geometric interpretation of $\mu_{p}(F, x)$. An application of Crofton's formula 7.11 shows, after a translation, that there is a small ball $B\left(x, r_{0}\right)$ on which

$$
\begin{align*}
\left(d d^{c} \log |F(z)-y|\right)^{p} \wedge\left(d d^{c} \log |z-x|\right)^{n-p} & = \\
\int_{S \in G(p, n)}\left(d d^{c} \log |F(z)-y|\right)^{p} & \wedge[x+S] d v(S) . \tag{9.7a}
\end{align*}
$$

For a rigorous proof of (9.7 a), we replace $\log |F(z)-y|$ by the "regularized" function $\frac{1}{2} \log \left(|F(z)-y|^{2}+\varepsilon^{2}\right)$ and let $\varepsilon$ tend to 0 on both sides. By (4.3) (resp. by (4.10)), the wedge product $\left(d d^{c} \log |F(z)-y|\right)^{p} \wedge[x+S]$ is well defined on a small ball $B\left(x, r_{0}\right)$ as soon as $x+S$ does not intersect $F^{-1}(y) \cap \partial B\left(x, r_{0}\right)$ (resp. intersects $F^{-1}(y) \cap B\left(x, r_{0}\right)$ at finitely many points); thanks to the assumption $\operatorname{codim}\left(F^{-1}(y), x\right) \geqslant p$, Sard's theorem shows that this is the case for all $S$ outside a negligible closed subset $E$ in $G(p, n)$ (resp. by Bertini, an analytic subset $A$ in $G(p, n)$ with $A \subset E)$. Fatou's lemma then implies that the inequality $\geqslant$ holds in (9.7 a). To get equality, we observe that we have bounded convergence on all complements $G(p, n) \backslash V(E)$ of neighborhoods $V(E)$ of $E$. However the mass of $\int_{V(E)}[x+S] d v(S)$ in $B\left(x, r_{0}\right)$ is proportional to $v(V(E))$ and therefore tends to 0 when $V(E)$ is small; this is sufficient to complete the proof, since Prop. 4.6 b ) gives

$$
\int_{z \in \bar{B}\left(x, r_{0}\right)}\left(d d^{c} \log \left(|F(z)-y|^{2}+\varepsilon^{2}\right)\right)^{p} \wedge \int_{S \in V(E)}[x+S] d v(S) \leqslant C v(V(E))
$$

with a constant $C$ independent of $\varepsilon$. By evaluating (9.7 a) on $\{x\}$, we get

$$
\begin{equation*}
\mu_{p}(F, x)=\int_{S \in G(p, n) \backslash A} \nu\left(\left(d d^{c} \log \left|F_{\lceil x+S}-z\right|\right)^{p}, x\right) d v(S) . \tag{9.7b}
\end{equation*}
$$

Let us choose a linear parametrization $g_{S}: \mathbb{C}^{p} \rightarrow S$ depending analytically on local coordinates of $S$ in $G(p, n)$. Then Theorem 8.4 with $T=\left[\mathbb{C}^{p}\right]$ and $\varphi(z, S)=\log \mid F \circ$ $g_{S}(z)-y \mid$ shows that

$$
\nu\left(\left(d d^{c} \log \left|F_{\lceil x+S}-z\right|\right)^{p}, x\right)=\nu\left(\left[\mathbb{C}^{p}\right], \log \left|F \circ g_{S}(z)-y\right|\right)
$$

is Zariski upper semicontinuous in $S$ on $G(p, n) \backslash A$. However, (9.6) shows that these numbers are integers, so $S \mapsto \nu\left(\left(d d^{c} \log \left|F_{\lceil x+S}-z\right|\right)^{p}, x\right)$ must be constant on a Zariski open subset in $G(p, n)$. By ( 9.7 b ), we obtain

$$
\begin{equation*}
\mu_{p}(F, x)=\mu_{p}\left(F_{\uparrow x+S}, x\right)=\text { local degree of } w^{\prime} \circ F_{\lceil x+S} \text { at } x \tag{9.7c}
\end{equation*}
$$

for generic subspaces $S \in G(p, n)$ and generic coordinates $w^{\prime}=\left(w_{1}, \ldots, w_{p}\right)$, $w^{\prime \prime}=$ $\left(w_{p+1}, \ldots, w_{m}\right)$ on $(Y, y)$.
(9.8) Example. Let $F: \mathbb{C}^{n} \longrightarrow \mathbb{C}^{n}$ be defined by

$$
F\left(z_{1}, \ldots, z_{n}\right)=\left(z_{1}^{s_{1}}, \ldots, z_{n}^{s_{n}}\right), \quad s_{1} \leqslant \ldots \leqslant s_{n}
$$

We claim that $\mu_{p}(F, 0)=s_{1} \ldots s_{p}$. In fact, for a generic $p$-dimensional subspace $S \subset \mathbb{C}^{n}$ such that $z_{1}, \ldots, z_{p}$ are coordinates on $S$ and $z_{p+1}, \ldots, z_{n}$ are linear forms in $z_{1}, \ldots, z_{p}$, and for generic coordinates $w^{\prime}=\left(w_{1}, \ldots, w_{p}\right)$, $w^{\prime \prime}=\left(w_{p+1}, \ldots, w_{n}\right)$ on $\mathbb{C}^{n}$, we can rearrange $w^{\prime}$ by linear combinations so that $w_{j} \circ F_{\upharpoonright S}$ is a linear combination of $\left(z_{j}^{s_{j}}, \ldots, z_{n}^{s_{n}}\right)$ and has non zero coefficient in $z_{j}^{s_{j}}$ as a polynomial in $\left(z_{j}, \ldots, z_{p}\right)$. It is then an exercise to show that $w^{\prime} \circ F_{\mid S}$ has covering degree $s_{1} \ldots s_{p}$ at 0 [compute inductively the roots $z_{n}, z_{n-1}, \ldots, z_{j}$ of $w_{j} \circ F_{\upharpoonright S}(z)=a_{j}$ and use Lemma II.3.10 to show that the $s_{j}$ values of $z_{j}$ lie near 0 when $\left(a_{1}, \ldots, a_{p}\right)$ are small].

We are now ready to prove the main result of this section, which describes the behaviour of Lelong numbers under proper morphisms. A similar weaker result was already proved in [Demailly 1982b] with some other non optimal multiplicities $\mu_{p}(F, x)$.
(9.9) Theorem. Let $T$ be a closed positive current of bidimension $(p, p)$ on $X$ and let $F: X \longrightarrow Y$ be an analytic map such that the restriction $F_{\uparrow \text { Supp } T}$ is proper. Let $I(y)$ be the set of points $x \in \operatorname{Supp} T \cap F^{-1}(y)$ such that $x$ is equal to its connected component in Supp $T \cap F^{-1}(y)$ and $\operatorname{codim}\left(F^{-1}(y), x\right) \geqslant p$. Then we have

$$
\nu\left(F_{\star} T, y\right) \geqslant \sum_{x \in I(y)} \mu_{p}(F, x) \nu(T, x)
$$

In particular, we have $\nu\left(F_{\star} T, y\right) \geqslant \sum_{x \in I(y)} \nu(T, x)$. This inequality no longer holds if the summation is extended to all points $x \in \operatorname{Supp} T \cap F^{-1}(Y)$ and if this set contains positive dimensional connected components: for example, if $F: X \longrightarrow Y$ contracts some exceptional subspace $E$ in $X$ to a point $y_{0}$ (e.g. if $F$ is a blow-up map, see $\S 7.12$ ), then $T=[E]$ has direct image $F_{\star}[E]=0$ thanks to (9.2).

Proof. We proceed in three steps.

Step 1. Reduction to the case of a single point $x$ in the fiber. It is sufficient to prove the inequality when the summation is taken over an arbitrary finite subset $\left\{x_{1}, \ldots, x_{N}\right\}$ of $I(y)$. As $x_{j}$ is equal to its connected component in Supp $T \cap F^{-1}(y)$, it has a fondamental system of relative open-closed neighborhoods, hence there are disjoint neighborhoods $U_{j}$ of $x_{j}$ such that $\partial U_{j}$ does not intersect Supp $T \cap F^{-1}(y)$. Then the image $F\left(\partial U_{j} \cap\right.$ Supp $\left.T\right)$ is a closed set which does not contain $y$. Let $W$ be a neighborhood of $y$ disjoint from all sets $F\left(\partial U_{j} \cap \operatorname{Supp} T\right)$, and let $V_{j}=U_{j} \cap F^{-1}(W)$. It is clear that $V_{j}$ is a neighborhood of $x_{j}$ and that $F_{\uparrow V_{j}}: V_{j} \rightarrow W$ has a proper restriction to Supp $T \cap V_{j}$. Moreover, we obviously have $F_{\star} T \geqslant \sum_{j}\left(F_{\left\lceil V_{j}\right.}\right)_{\star} T$ on $W$. Therefore, it is enough to check the inequality $\nu\left(F_{\star} T, y\right) \geqslant \mu_{p}(F, x) \nu(T, x)$ for a single point $x \in I(y)$, in the case when $X \subset \mathbb{C}^{n}, Y \subset \mathbb{C}^{m}$ are open subsets and $x=y=0$.

Step 2. Reduction to the case when $F$ is finite. By (9.4), we have

$$
\begin{aligned}
\nu\left(F_{\star} T, 0\right) & =\inf _{V \ni 0} \int_{V} T \wedge\left(d d^{c} \log |F|\right)^{p} \\
& =\inf _{V \ni 0} \lim _{\varepsilon \rightarrow 0} \int_{V} T \wedge\left(d d^{c} \log \left(|F|+\varepsilon|z|^{N}\right)\right)^{p},
\end{aligned}
$$

and the integrals are well defined as soon as $\partial V$ does not intersect the set Supp $T \cap F^{-1}(0)$ (may be after replacing $\log |F|$ by $\max \{\log |F|, s\}$ with $s \ll 0$ ). For every $V$ and $\varepsilon$, the last integral is larger than $\nu\left(G_{\star} T, 0\right)$ where $G$ is the finite morphism defined by

$$
G: X \longrightarrow Y \times \mathbb{C}^{n}, \quad\left(z_{1}, \ldots, z_{n}\right) \longmapsto\left(F_{1}(z), \ldots, F_{m}(z), z_{1}^{N}, \ldots, z_{n}^{N}\right)
$$

We claim that for $N$ large enough we have $\mu_{p}(F, 0)=\mu_{p}(G, 0)$. In fact, $x \in I(y)$ implies by definition $\operatorname{codim}\left(F^{-1}(0), 0\right) \geqslant p$. Hence, if $S=\left\{u_{1}=\ldots=u_{n-p}=0\right\}$ is a generic $p$-dimensional subspace of $\mathbb{C}^{n}$, the germ of variety $F^{-1}(0) \cap S$ defined by $\left(F_{1}, \ldots, F_{m}, u_{1}, \ldots, u_{n-p}\right)$ is $\{0\}$. Hilbert's Nullstellensatz implies that some powers of $z_{1}, \ldots, z_{n}$ are in the ideal $\left(F_{j}, u_{k}\right)$. Therefore $|F(z)|+|u(z)| \geqslant C|z|^{a}$ near 0 for some integer $a$ independent of $S$ (to see this, take coefficients of the $u_{k}$ 's as additional variables); in particular $|F(z)| \geqslant C|z|^{a}$ for $z \in S$ near 0 . The comparison theorem 7.1 then shows that $\mu_{p}(F, 0)=\mu_{p}(G, 0)$ for $N \geqslant a$. If we are able to prove that $\nu\left(G_{\star} T, 0\right) \geqslant$ $\mu_{p}(G, 0) \nu(T, 0)$ in case $G$ is finite, the obvious inequality $\nu\left(F_{\star} T, 0\right) \geqslant \nu\left(G_{\star} T, 0\right)$ concludes the proof.
Step 3. Proof of the inequality $\nu\left(F_{\star} T, y\right) \geqslant \mu_{p}(F, x) \nu(T, x)$ when $F$ is finite and $F^{-1}(y)=$ $x$. Then $\varphi(z)=\log |F(z)-y|$ has a single isolated pole at $x$ and we have $\mu_{p}(F, x)=$ $\nu\left(\left(d d^{c} \varphi\right)^{p}, x\right)$. It is therefore sufficient to apply to following Proposition.
(9.10) Proposition. Let $\varphi$ be a semi-exhaustive continuous plurisubharmonic function on $X$ with a single isolated pole at $x$. Then

$$
\nu(T, \varphi) \geqslant \nu(T, x) \nu\left(\left(d d^{c} \varphi\right)^{p}, x\right) .
$$

Proof. Since the question is local, we can suppose that $X$ is the ball $B\left(0, r_{0}\right)$ in $\mathbb{C}^{n}$ and $x=0$. Set $X^{\prime}=B\left(0, r_{1}\right)$ with $r_{1}<r_{0}$ and $\Phi(z, g)=\varphi \circ g(z)$ for $g \in \mathrm{Gl}_{n}(\mathbb{C})$. Then there is a small neighborhood $\Omega$ of the unitary group $U(n) \subset \mathrm{Gl}_{n}(\mathbb{C})$ such that $\Phi$ is plurisubharmonic on $X^{\prime} \times \Omega$ and semi-exhaustive with respect to $X^{\prime}$. Theorem 8.4
implies that the map $g \mapsto \nu(T, \varphi \circ g)$ is Zariski upper semi-continuous on $\Omega$. In particular, we must have $\nu(T, \varphi \circ g) \leqslant \nu(T, \varphi)$ for all $g \in \Omega \backslash A$ in the complement of a complex analytic set $A$. Since $\mathrm{Gl}_{n}(\mathbb{C})$ is the complexification of $U(n)$, the intersection $U(n) \cap A$ must be a nowhere dense real analytic subset of $U(n)$. Therefore, if $d v$ is the Haar measure of mass 1 on $U(n)$, we have

$$
\begin{align*}
\nu(T, \varphi) & \geqslant \int_{g \in U(n)} \nu(T, \varphi \circ g) d v(g) \\
& =\lim _{r \rightarrow 0} \int_{g \in U(n)} d v(g) \int_{B(0, r)} T \wedge\left(d d^{c} \varphi \circ g\right)^{p} . \tag{9.11}
\end{align*}
$$

Since $\int_{g \in U(n)}\left(d d^{c} \varphi \circ g\right)^{p} d v(g)$ is a unitary invariant $(p, p)$-form on $B$, Lemma 7.10 implies

$$
\int_{g \in U(n)}\left(d d^{c} \varphi \circ g\right)^{p} d v(g)=\left(d d^{c} \chi(\log |z|)\right)^{p}
$$

where $\chi$ is a convex increasing function. The Lelong number at 0 of the left hand side is equal to $\nu\left(\left(d d^{c} \varphi\right)^{p}, 0\right)$, and must be equal to the Lelong number of the right hand side, which is $\lim _{t \rightarrow-\infty} \chi^{\prime}(t)^{p}$ (to see this, use either Formula (5.5) or Th. 7.8). Thanks to the last equality, Formulas (9.11) and (5.5) imply

$$
\begin{aligned}
\nu(T, \varphi) & \geqslant \lim _{r \rightarrow 0} \int_{B(0, r)} T \wedge\left(d d^{c} \chi(\log |z|)\right)^{p} \\
& =\lim _{r \rightarrow 0} \chi^{\prime}(\log r-0)^{p} \nu(T, 0, r) \geqslant \nu\left(\left(d d^{c} \varphi\right)^{p}, 0\right) \nu(T, 0)
\end{aligned}
$$

Another interesting question is to know whether it is possible to get inequalities in the opposite direction, i.e. to find upper bounds for $\nu\left(F_{\star} T, y\right)$ in terms of the Lelong numbers $\nu(T, x)$. The example $T=[\Gamma]$ with the curve $\Gamma: t \mapsto\left(t^{a}, t^{a+1}, t\right)$ in $\mathbb{C}^{3}$ and $F: \mathbb{C}^{3} \rightarrow \mathbb{C}^{2},\left(z_{1}, z_{2}, z_{3}\right) \mapsto\left(z_{1}, z_{2}\right)$, for which $\nu(T, 0)=1$ and $\nu\left(F_{\star} T, 0\right)=a$, shows that this may be possible only when $F$ is finite. In this case, we have:
(9.12) Theorem. Let $F: X \rightarrow Y$ be a proper and finite analytic map and let $T$ be $a$ closed positive current of bidimension $(p, p)$ on $X$. Then
(a)

$$
\nu\left(F_{\star} T, y\right) \leqslant \sum_{x \in \operatorname{Supp} T \cap F^{-1}(y)} \bar{\mu}_{p}(F, x) \nu(T, x)
$$

where $\bar{\mu}_{p}(F, x)$ is the multiplicity defined as follows: if $H:(X, x) \rightarrow\left(\mathbb{C}^{n}, 0\right)$ is a germ of finite map, we set

$$
\begin{align*}
\sigma(H, x) & =\inf \left\{\alpha>0 ; \exists C>0,|H(z)| \geqslant C|z-x|^{\alpha} \text { near } x\right\},  \tag{b}\\
\bar{\mu}_{p}(F, x) & =\inf _{G} \frac{\sigma(G \circ F, x)^{p}}{\mu_{p}(G, 0)}, \tag{c}
\end{align*}
$$

where $G$ runs over all germs of maps $(Y, y) \longrightarrow\left(\mathbb{C}^{n}, 0\right)$ such that $G \circ F$ is finite.
Proof. If $F^{-1}(y)=\left\{x_{1}, \ldots, x_{N}\right\}$, there is a neighborhood $W$ of $y$ and disjoint neighborhoods $V_{j}$ of $x_{j}$ such that $F^{-1}(W)=\bigcup V_{j}$. Then $F_{\star} T=\sum\left(F_{\mid V_{j}}\right)_{\star} T$ on $W$, so it
is enough to consider the case when $F^{-1}(y)$ consists of a single point $x$. Therefore, we assume that $F: V \rightarrow W$ is proper and finite, where $V, W$ are neighborhoods of 0 in $\mathbb{C}^{n}$, $\mathbb{C}^{m}$ and $F^{-1}(0)=\{0\}$. Let $G:\left(\mathbb{C}^{m}, 0\right) \longrightarrow\left(\mathbb{C}^{n}, 0\right)$ be a germ of map such that $G \circ F$ is finite. Hilbert's Nullstellensatz shows that there exists $\alpha>0$ and $C>0$ such that $|G \circ F(z)| \geqslant C|z|^{\alpha}$ near 0 . Then the comparison theorem 7.1 implies

$$
\nu\left(G_{\star} F_{\star} T, 0\right)=\nu(T, \log |G \circ F|) \leqslant \alpha^{p} \nu(T, \log |z|)=\alpha^{p} \nu(T, 0)
$$

On the other hand, Th. 9.9 applied to $\Theta=F_{\star} T$ on $W$ gives

$$
\nu\left(G_{\star} F_{\star} T, 0\right) \geqslant \mu_{p}(G, 0) \nu\left(F_{\star} T, 0\right)
$$

Therefore

$$
\nu\left(F_{\star} T, 0\right) \leqslant \frac{\alpha^{p}}{\mu_{p}(G, 0)} \nu(T, 0) .
$$

The infimum of all possible values of $\alpha$ is by definition $\sigma(G \circ F, 0)$, thus by taking the infimum over $G$ we obtain

$$
\nu\left(F_{\star} T, 0\right) \leqslant \bar{\mu}_{p}(F, 0) \nu(T, 0)
$$

(9.13) Example. Let $F\left(z_{1}, \ldots, z_{n}\right)=\left(z_{1}^{s_{1}}, \ldots, z_{n}^{s_{n}}\right), s_{1} \leqslant \ldots \leqslant s_{n}$ as in 9.8. Then we have

$$
\mu_{p}(F, 0)=s_{1} \ldots s_{p}, \quad \bar{\mu}_{p}(F, 0)=s_{n-p+1} \ldots s_{n}
$$

To see this, let $s$ be the lowest common multiple of $s_{1}, \ldots, s_{n}$ and let $G\left(z_{1}, \ldots, z_{n}\right)=$ $\left(z_{1}^{s / s_{1}}, \ldots, z_{n}^{s / s_{n}}\right)$. Clearly $\mu_{p}(G, 0)=\left(s / s_{n-p+1}\right) \ldots\left(s / s_{n}\right)$ and $\sigma(G \circ F, 0)=s$, so we get by definition $\bar{\mu}_{p}(F, 0) \leqslant s_{n-p+1} \ldots s_{n}$. Finally, if $T=[A]$ is the current of integration over the $p$-dimensional subspace $A=\left\{z_{1}=\ldots=z_{n-p}=0\right\}$, then $F_{\star}[A]=s_{n-p+1} \ldots s_{n}[A]$ because $F_{\upharpoonright A}$ has covering degree $s_{n-p+1} \ldots s_{n}$. Theorem 9.12 shows that we must have $s_{n-p+1} \ldots s_{n} \leqslant \bar{\mu}_{p}(F, 0)$, QED. If $\lambda_{1} \leqslant \ldots \leqslant \lambda_{n}$ are positive real numbers and $s_{j}$ is taken to be the integer part of $k \lambda_{j}$ as $k$ tends to $+\infty$, Theorems 9.9 and 9.12 imply in the limit the following:
(9.14) Corollary. For $0<\lambda_{1} \leqslant \ldots \leqslant \lambda_{n}$, Kiselman's directional Lelong numbers satisfy the inequalities

$$
\lambda_{1} \ldots \lambda_{p} \nu(T, x) \leqslant \nu(T, x, \lambda) \leqslant \lambda_{n-p+1} \ldots \lambda_{n} \nu(T, x)
$$

(9.15) Remark. It would be interesting to have a direct geometric interpretation of $\bar{\mu}_{p}(F, x)$. In fact, we do not even know whether $\bar{\mu}_{p}(F, x)$ is always an integer.

## § 10. A Schwarz Lemma. Application to Number Theory

In this section, we show how Jensen's formula and Lelong numbers can be used to prove a fairly general Schwarz lemma relating growth and zeros of entire functions in $\mathbb{C}^{n}$. In order to simplify notations, we denote by $|F|_{r}$ the supremum of the modulus of a function $F$ on the ball of center 0 and radius $r$. Then, following [Demailly 1982a], we present some applications with a more arithmetical flavour.
(10.1) Schwarz lemma. Let $P_{1}, \ldots, P_{N} \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$ be polynomials of degree $\delta$, such that their homogeneous parts of degree $\delta$ do not vanish simultaneously except at 0 . Then there is a constant $C \geqslant 2$ such that for all entire functions $F \in \mathscr{O}\left(\mathbb{C}^{n}\right)$ and all $R \geqslant r \geqslant 1$ we have

$$
\log |F|_{r} \leqslant \log |F|_{R}-\delta^{1-n} \nu\left(\left[Z_{F}\right], \log |P|\right) \log \frac{R}{C r}
$$

where $Z_{F}$ is the zero divisor of $F$ and $P=\left(P_{1}, \ldots, P_{N}\right): \mathbb{C}^{n} \longrightarrow \mathbb{C}^{N}$. Moreover

$$
\nu\left(\left[Z_{F}\right], \log |P|\right) \geqslant \sum_{w \in P^{-1}(0)} \operatorname{ord}(F, w) \mu_{n-1}(P, w)
$$

where $\operatorname{ord}(F, w)$ denotes the vanishing order of $F$ at $w$ and $\mu_{n-1}(P, w)$ is the $(n-1)$ multiplicity of $P$ at $w$, as defined in (9.5) and (9.7).

Proof. Our assumptions imply that $P$ is a proper and finite map. The last inequality is then just a formal consequence of formula (9.4) and Th. 9.9 applied to $T=\left[Z_{F}\right]$. Let $Q_{j}$ be the homogeneous part of degree $\delta$ in $P_{j}$. For $z_{0} \in B(0, r)$, we introduce the weight functions

$$
\varphi(z)=\log |P(z)|, \quad \psi(z)=\log \left|Q\left(z-z_{0}\right)\right| .
$$

Since $Q^{-1}(0)=\{0\}$ by hypothesis, the homogeneity of $Q$ shows that there are constants $C_{1}, C_{2}>0$ such that

$$
\begin{equation*}
C_{1}|z|^{\delta} \leqslant|Q(z)| \leqslant C_{2}|z|^{\delta} \quad \text { on } \mathbb{C}^{n} \tag{10.2}
\end{equation*}
$$

The homogeneity also implies $\left(d d^{c} \psi\right)^{n}=\delta^{n} \delta_{z_{0}}$. We apply the Lelong Jensen formula 6.5 to the measures $\mu_{\psi, s}$ associated with $\psi$ and to $V=\log |F|$. This gives

$$
\begin{equation*}
\mu_{\psi, s}(\log |F|)-\delta^{n} \log \left|F\left(z_{0}\right)\right|=\int_{-\infty}^{s} d t \int_{\{\psi<t\}}\left[Z_{F}\right] \wedge\left(d d^{c} \psi\right)^{n-1} \tag{10.3}
\end{equation*}
$$

By (6.2), $\mu_{\psi, s}$ has total mass $\delta^{n}$ and has support in

$$
\{\psi(z)=s\}=\left\{Q\left(z-z_{0}\right)=e^{s}\right\} \subset B\left(0, r+\left(e^{s} / C_{1}\right)^{1 / \delta}\right)
$$

Note that the inequality in the Schwarz lemma is obvious if $R \leqslant C r$, so we can assume $R \geqslant C r \geqslant 2 r$. We take $s=\delta \log (R / 2)+\log C_{1}$; then

$$
\{\psi(z)=s\} \subset B(0, r+R / 2) \subset B(0, R)
$$

In particular, we get $\mu_{\psi, s}(\log |F|) \leqslant \delta^{n} \log |F|_{R}$ and formula (10.3) gives

$$
\begin{equation*}
\log |F|_{R}-\log \left|F\left(z_{0}\right)\right| \geqslant \delta^{-n} \int_{s_{0}}^{s} d t \int_{\{\psi<t\}}\left[Z_{F}\right] \wedge\left(d d^{c} \psi\right)^{n-1} \tag{10.4}
\end{equation*}
$$

for any real number $s_{0}<s$. The proof will be complete if we are able to compare the integral in (10.4) to the corresponding integral with $\varphi$ in place of $\psi$. The argument for this is quite similar to the proof of the comparison theorem, if we observe that $\psi \sim \varphi$ at infinity. We introduce the auxiliary function

$$
w= \begin{cases}\max \{\psi,(1-\varepsilon) \varphi+\varepsilon t-\varepsilon\} & \text { on }\{\psi \geqslant t-2\}, \\ (1-\varepsilon) \varphi+\varepsilon t-\varepsilon & \text { on }\{\psi \leqslant t-2\},\end{cases}
$$

with a constant $\varepsilon$ to be determined later, such that $(1-\varepsilon) \varphi+\varepsilon t-\varepsilon>\psi$ near $\{\psi=t-2\}$ and $(1-\varepsilon) \varphi+\varepsilon t-\varepsilon<\psi$ near $\{\psi=t\}$. Then Stokes' theorem implies

$$
\begin{aligned}
& \int_{\{\psi<t\}}\left[Z_{F}\right] \wedge\left(d d^{c} \psi\right)^{n-1}=\int_{\{\psi<t\}}\left[Z_{F}\right] \wedge\left(d d^{c} w\right)^{n-1} \\
& \quad \geqslant(1-\varepsilon)^{n-1} \int_{\{\psi<t-2\}}\left[Z_{F}\right] \wedge\left(d d^{c} \varphi\right)^{n-1} \geqslant(1-\varepsilon)^{n-1} \nu\left(\left[Z_{F}\right], \log |P|\right)
\end{aligned}
$$

By (10.2) and our hypothesis $\left|z_{0}\right|<r$, the condition $\psi(z)=t$ implies

$$
\begin{aligned}
&\left|Q\left(z-z_{0}\right)\right|=e^{t} \Longrightarrow e^{t / \delta} / C_{1}^{1 / \delta} \leqslant\left|z-z_{0}\right| \leqslant e^{t / \delta} / C_{2}^{1 / \delta} \\
&\left|P(z)-Q\left(z-z_{0}\right)\right| \leqslant C_{3}\left(1+\left|z_{0}\right|\right)\left(1+|z|+\left|z_{0}\right|\right)^{\delta-1} \leqslant C_{4} r\left(r+e^{t / \delta}\right)^{\delta-1} \\
&\left|\frac{P(z)}{Q\left(z-z_{0}\right)}-1\right| \leqslant C_{4} r e^{-t / \delta}\left(r e^{-t / \delta}+1\right)^{\delta-1} \leqslant 2^{\delta-1} C_{4} r e^{-t / \delta}
\end{aligned}
$$

provided that $t \geqslant \delta \log r$. Hence for $\psi(z)=t \geqslant s_{0} \geqslant \delta \log \left(2^{\delta} C_{4} r\right)$, we get

$$
|\varphi(z)-\psi(z)|=\left|\log \frac{|P(z)|}{\left|Q\left(z-z_{0}\right)\right|}\right| \leqslant C_{5} r e^{-t / \delta}
$$

Now, we have

$$
[(1-\varepsilon) \varphi+\varepsilon t-\varepsilon]-\psi=(1-\varepsilon)(\varphi-\psi)+\varepsilon(t-1-\psi)
$$

so this difference is $<C_{5} r e^{-t / \delta}-\varepsilon$ on $\{\psi=t\}$ and $>-C_{5} r e^{(2-t) / \delta}+\varepsilon$ on $\{\psi=t-2\}$. Hence it is sufficient to take $\varepsilon=C_{5} r e^{(2-t) / \delta}$. This number has to be $<1$, so we take $t \geqslant s_{0} \geqslant 2+\delta \log \left(C_{5} r\right)$. Moreover, (10.5) actually holds only if $P^{-1}(0) \subset\{\psi<t-2\}$, so by (10.2) it is enough to take $t \geqslant s_{0} \geqslant 2+\log \left(C_{2}\left(r+C_{6}\right)^{\delta}\right)$ where $C_{6}$ is such that $P^{-1}(0) \subset \bar{B}\left(0, C_{6}\right)$. Finally, we see that we can choose

$$
s=\delta \log R-C_{7}, \quad s_{0}=\delta \log r+C_{8}
$$

and inequalities (10.4), (10.5) together imply

$$
\log |F|_{R}-\log \left|F\left(z_{0}\right)\right| \geqslant \delta^{-n}\left(\int_{s_{0}}^{s}\left(1-C_{5} r e^{(2-t) / \delta}\right)^{n-1} d t\right) \nu\left(\left[Z_{F}\right], \log |P|\right) .
$$

The integral is bounded below by

$$
\int_{C_{8}}^{\delta \log (R / r)-C_{7}}\left(1-C_{9} e^{-t / \delta}\right) d t \geqslant \delta \log (R / C r)
$$

This concludes the proof, by taking the infimum when $z_{0}$ runs over $B(0, r)$.
(10.6) Corollary. Let $S$ be a finite subset of $\mathbb{C}^{n}$ and let $\delta$ be the minimal degree of algebraic hypersurfaces containing $S$. Then there is a constant $C \geqslant 2$ such that for all $F \in \mathbb{O}\left(\mathbb{C}^{n}\right)$ and all $R \geqslant r \geqslant 1$ we have

$$
\log |F|_{r} \leqslant \log |F|_{R}-\operatorname{ord}(F, S) \frac{\delta+n(n-1) / 2}{n!} \log \frac{R}{C r}
$$

where $\operatorname{ord}(F, S)=\min _{w \in S} \operatorname{ord}(F, w)$.
Proof. In view of Th. 10.1, we only have to select suitable polynomials $P_{1}, \ldots, P_{N}$. The vector space $\mathbb{C}\left[z_{1}, \ldots, z_{n}\right]_{<\delta}$ of polynomials of degree $<\delta$ in $\mathbb{C}^{n}$ has dimension

$$
m(\delta)=\binom{\delta+n-1}{n}=\frac{\delta(\delta+1) \ldots(\delta+n-1)}{n!}
$$

By definition of $\delta$, the linear forms

$$
\mathbb{C}\left[z_{1}, \ldots, z_{n}\right]_{<\delta} \longrightarrow \mathbb{C}, \quad P \longmapsto P(w), \quad w \in S
$$

vanish simultaneously only when $P=0$. Hence we can find $m=m(\delta)$ points $w_{1}, \ldots$, $w_{m} \in S$ such that the linear forms $P \mapsto P\left(w_{j}\right)$ define a basis of $\mathbb{C}\left[z_{1}, \ldots, z_{n}\right]_{<\delta \delta}^{\star}$. This means that there is a unique polynomial $P \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right]_{<\delta}$ which takes given values $P\left(w_{j}\right)$ for $1 \leqslant j \leqslant m$. In particular, for every multiindex $\alpha,|\alpha|=\delta$, there is a unique polynomial $R_{\alpha} \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right]_{<\delta}$ such that $R_{\alpha}\left(w_{j}\right)=w_{j}^{\alpha}$. Then the polynomials $P_{\alpha}(z)=z^{\alpha}-R_{\alpha}(z)$ have degree $\delta$, vanish at all points $w_{j}$ and their homogeneous parts of maximum degree $Q_{\alpha}(z)=z^{\alpha}$ do not vanish simultaneously except at 0 . We simply use the fact that $\mu_{n-1}\left(P, w_{j}\right) \geqslant 1$ to get

$$
\nu\left(\left[Z_{F}\right], \log |P|\right) \geqslant \sum_{w \in P^{-1}(0)} \operatorname{ord}(F, w) \geqslant m(\delta) \operatorname{ord}(F, S) .
$$

Theorem 10.1 then gives the desired inequality, because $m(\delta)$ is a polynomial with positive coefficients and with leading terms

$$
\frac{1}{n!}\left(\delta^{n}+n(n-1) / 2 \delta^{n-1}+\ldots\right)
$$

Let $S$ be a finite subset of $\mathbb{C}^{n}$. According to [Waldschmidt 1976], we introduce for every integer $t>0$ a number $\omega_{t}(S)$ equal to the minimal degree of polynomials $P \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$ which vanish at order $\geqslant t$ at every point of $S$. The obvious subadditivity property

$$
\omega_{t_{1}+t_{2}}(S) \leqslant \omega_{t_{1}}(S)+\omega_{t_{2}}(S)
$$

easily shows that

$$
\Omega(S):=\inf _{t>0} \frac{\omega_{t}(S)}{t}=\lim _{t \rightarrow+\infty} \frac{\omega_{t}(S)}{t}
$$

We call $\omega_{1}(S)$ the degree of $S$ (minimal degree of algebraic hypersurfaces containing $S$ ) and $\Omega(S)$ the singular degree of $S$. If we apply Cor. 10.6 to a polynomial $F$ vanishing at order $t$ on $S$ and fix $r=1$, we get

$$
\log |F|_{R} \geqslant t \frac{\delta+n(n-1) / 2}{n!} \log \frac{R}{C}+\log |F|_{1}
$$

with $\delta=\omega_{1}(S)$, in particular

$$
\operatorname{deg} F \geqslant t \frac{\omega_{1}(S)+n(n-1) / 2}{n!}
$$

The minimum of $\operatorname{deg} F$ over all such $F$ is by definition $\omega_{t}(S)$. If we divide by $t$ and take the infimum over $t$, we get the interesting inequality

$$
\begin{equation*}
\frac{\omega_{t}(S)}{t} \geqslant \Omega(S) \geqslant \frac{\omega_{1}(S)+n(n-1) / 2}{n!} \tag{10.7}
\end{equation*}
$$

(10.8) Remark. The constant $\frac{\omega_{1}(S)+n(n-1) / 2}{n!}$ in (10.6) and (10.7) is optimal for $n=1,2$ but not for $n \geqslant 3$. It can be shown by means of Hörmander's $L^{2}$ estimates [Waldschmidt 1978] that for every $\varepsilon>0$ the Schwarz lemma (10.6) holds with coefficient $\Omega(S)-\varepsilon$ :

$$
\log |F|_{r} \leqslant \log |F|_{R}-\operatorname{ord}(F, S)(\Omega(S)-\varepsilon) \log \frac{R}{C_{\varepsilon} r}
$$

and that $\Omega(S) \geqslant\left(\omega_{u}(S)+1\right) /(u+n-1)$ for every $u \geqslant 1$; this last inequality is due to [Esnault-Viehweg 1983], who used deep tools of algebraic geometry; [Azhari 1990] reproved it recently by means of Hörmander's $L^{2}$ estimates. Rather simple examples [Demailly 1982a] lead to the conjecture

$$
\Omega(S) \geqslant \frac{\omega_{u}(S)+n-1}{u+n-1} \quad \text { for every } \quad u \geqslant 1
$$

The special case $u=1$ of the conjecture was first stated by [Chudnovsky 1979].

Finally, let us mention that Cor. 10.6 contains Bombieri's theorem on algebraic values of meromorphic maps satisfying algebraic differential equations [Bombieri 1970]. Recall that an entire function $F \in \mathbb{O}\left(\mathbb{C}^{n}\right)$ is said to be of order $\leqslant \rho$ if for every $\varepsilon>0$ there is a constant $C_{\varepsilon}$ such that $|F(z)| \leqslant C_{\varepsilon} \exp \left(|z|^{\rho+\varepsilon}\right)$. A meromorphic function is said to be of order $\leqslant \rho$ if it can be written $G / H$ where $G, H$ are entire functions of order $\leqslant \rho$.
(10.9) Theorem ([Bombieri 1970]). Let $F_{1}, \ldots, F_{N}$ be meromorphic functions on $\mathbb{C}^{n}$, such that $F_{1}, \ldots, F_{d}, n<d \leqslant N$, are algebraically independent over $\mathbb{Q}$ and have $f_{i}$ nite orders $\rho_{1}, \ldots, \rho_{d}$. Let $K$ be a number field of degree $[K: \mathbb{Q}]$. Suppose that the ring $K\left[f_{1}, \ldots, f_{N}\right]$ is stable under all derivations $d / d z_{1}, \ldots, d / d z_{n}$. Then the set $S$ of
points $z \in \mathbb{C}^{n}$, distinct from the poles of the $F_{j}$ 's, such that $\left(F_{1}(z), \ldots, F_{N}(z)\right) \in K^{N}$ is contained in an algebraic hypersurface whose degree $\delta$ satisfies

$$
\frac{\delta+n(n-1) / 2}{n!} \leqslant \frac{\rho_{1}+\ldots+\rho_{d}}{d-n}[K: \mathbb{Q}] .
$$

Proof. If the set $S$ is not contained in any algebraic hypersurface of degree $<\delta$, the linear algebra argument used in the proof of Cor. 10.6 shows that we can find $m=m(\delta)$ points $w_{1}, \ldots, w_{m} \in S$ which are not located on any algebraic hypersurface of degree $<\delta$. Let $H_{1}, \ldots, H_{d}$ be the denominators of $F_{1}, \ldots, F_{d}$. The standard arithmetical methods of transcendental number theory allow us to construct a sequence of entire functions in the following way: we set

$$
G=P\left(F_{1}, \ldots, F_{d}\right)\left(H_{1} \ldots H_{d}\right)^{s}
$$

where $P$ is a polynomial of degree $\leqslant s$ in each variable with integer coefficients. The polynomials $P$ are chosen so that $G$ vanishes at a very high order at each point $w_{j}$. This amounts to solving a linear system whose unknowns are the coefficients of $P$ and whose coefficients are polynomials in the derivatives of the $F_{j}$ 's (hence lying in the number field $K)$. Careful estimates of size and denominators and a use of the Dirichlet-Siegel box principle lead to the following lemma, see e.g. [Waldschmidt 1978].
(10.10) Lemma. For every $\varepsilon>0$, there exist constants $C_{1}, C_{2}>0, r \geqslant 1$ and an infinite sequence $G_{t}$ of entire functions, $t \in T \subset \mathbb{N}$ (depending on $m$ and on the choice of the points $\left.w_{j}\right)$, such that
a) $G_{t}$ vanishes at order $\geqslant t$ at all points $w_{1}, \ldots, w_{m}$;
b) $\left|G_{t}\right|_{r} \geqslant\left(C_{1} t\right)^{-t[K: \mathbb{Q}]}$;
c) $\left|G_{t}\right|_{R(t)} \leqslant C_{2}^{t} \quad$ where $R(t)=\left(t^{d-n} / \log t\right)^{1 /\left(\rho_{1}+\ldots+\rho_{d}+\varepsilon\right)}$.

An application of Cor. 10.6 to $F=G_{t}$ and $R=R(t)$ gives the desired bound for the degree $\delta$ as $t$ tends to $+\infty$ and $\varepsilon$ tends to 0 . If $\delta_{0}$ is the largest integer which satisfies the inequality of Th. 10.9, we get a contradiction if we take $\delta=\delta_{0}+1$. This shows that $S$ must be contained in an algebraic hypersurface of degree $\delta \leqslant \delta_{0}$.

## Chapter IV

## Sheaf Cohomology and Spectral Sequences


#### Abstract

One of the main topics of this book is the computation of various cohomology groups arising in algebraic geometry. The theory of sheaves provides a general framework in which many cohomology theories can be treated in a unified way. The cohomology theory of sheaves will be constructed here by means of Godement's simplicial flabby resolution. However, we have emphasized the analogy with Alexander-Spanier cochains in order to give a simple definition of the cup product. In this way, all the basic properties of cohomology groups (long exact sequences, Mayer Vietoris exact sequence, Leray's theorem, relations with Cech cohomology, De Rham-Weil isomorphism theorem) can be derived in a very elementary way from the definitions. Spectral sequences and hypercohomology groups are then introduced, with two principal examples in view: the Leray spectral sequence and the Hodge-Frölicher spectral sequence. The basic results concerning cohomology groups with constant or locally constant coefficients (invariance by homotopy, Poincaré duality, Leray-Hirsch theorem) are also included, in order to present a self-contained approach of algebraic topology.


## § 1. Basic Results of Homological Algebra

Let us first recall briefly some standard notations and results of homological algebra that will be used systematically in the sequel. Let $R$ be a commutative ring with unit. A differential module $(K, d)$ is a $R$-module $K$ together with an endomorphism $d: K \rightarrow K$, called the differential, such that $d \circ d=0$. The modules of cycles and of boundaries of $K$ are defined respectively by

$$
\begin{equation*}
Z(K)=\operatorname{ker} d, \quad B(K)=\operatorname{Im} d \tag{1.1}
\end{equation*}
$$

Our hypothesis $d \circ d=0$ implies $B(K) \subset Z(K)$. The homology group of $K$ is by definition the quotient module

$$
\begin{equation*}
H(K)=Z(K) / B(K) \tag{1.2}
\end{equation*}
$$

A morphism of differential modules $\varphi: K \longrightarrow L$ is a $R$-homomorphism $\varphi: K \longrightarrow L$ such that $d \circ \varphi=\varphi \circ d$; here we denote by the same symbol $d$ the differentials of $K$ and $L$. It is then clear that $\varphi(Z(K)) \subset Z(L)$ and $\varphi(B(K)) \subset B(L)$. Therefore, we get an induced morphism on homology groups, denoted

$$
\begin{equation*}
H(\varphi): H(K) \longrightarrow H(L) \tag{1.3}
\end{equation*}
$$

It is easily seen that $H$ is a functor, i.e. $H(\psi \circ \varphi)=H(\psi) \circ H(\varphi)$. We say that two morphisms $\varphi, \psi: K \longrightarrow L$ are homotopic if there exists a $R$-linear map $h: K \longrightarrow L$ such that

$$
\begin{equation*}
d \circ h+h \circ d=\psi-\varphi . \tag{1.4}
\end{equation*}
$$

Then $h$ is said to be a homotopy between $\varphi$ and $\psi$. For every cocycle $z \in Z(K)$, we infer $\psi(z)-\varphi(z)=d h(z)$, hence the maps $H(\varphi)$ and $H(\psi)$ coincide. The module $K$ itself is said to be homotopic to 0 if $\operatorname{Id}_{K}$ is homotopic to 0 ; then $H(K)=0$.
(1.5) Snake lemma. Let

$$
0 \longrightarrow K \xrightarrow{\varphi} L \xrightarrow{\psi} M \longrightarrow 0
$$

be a short exact sequence of morphisms of differential modules. Then there exists a homomorphism $\partial: H(M) \longrightarrow H(K)$, called the connecting homomorphism, and a homology exact sequence

$$
H(K) \xrightarrow{H(\varphi)} H(L) \xrightarrow{H(\psi)} H(M)
$$

Moreover, to any commutative diagram of short exact sequences

is associated a commutative diagram of homology exact sequences

$$
\begin{array}{cccc}
H(K) & \longrightarrow H(L) & \longrightarrow H(M) & \xrightarrow{\partial} H(K) \\
\downarrow & \downarrow & \downarrow & \downarrow \\
H(\widetilde{K}) \longrightarrow H(\widetilde{L}) & \longrightarrow H(\widetilde{M}) & \xrightarrow{\partial} H(\widetilde{K}) \longrightarrow \cdots
\end{array}
$$

Proof. We first define the connecting homomorphism $\partial$ : let $m \in Z(M)$ represent a given cohomology class $\{m\}$ in $H(M)$. Then

$$
\partial\{m\}=\{k\} \in H(K)
$$

is the class of any element $k \in \varphi^{-1} d \psi^{-1}(m)$, as obtained through the following construction:

$$
\begin{aligned}
& l \in L \stackrel{\psi}{\longmapsto} m \in M \\
& \downarrow d \quad \downarrow d \\
& k \in K \stackrel{\varphi}{\longmapsto} d l \in L \stackrel{\psi}{\longmapsto} 0 \in M .
\end{aligned}
$$

The element $l$ is chosen to be a preimage of $m$ by the surjective map $\psi$; as $\psi(d l)=$ $d(m)=0$, there exists a unique element $k \in K$ such that $\varphi(k)=d l$. The element $k$ is actually a cocycle in $Z(K)$ because $\varphi$ is injective and

$$
\varphi(d k)=d \varphi(k)=d(d l)=0 \quad \Longrightarrow \quad d k=0
$$

The map $\partial$ will be well defined if we show that the cohomology class $\{k\}$ depends only on $\{m\}$ and not on the choices made for the representatives $m$ and $l$. Consider another
representative $m^{\prime}=m+d m_{1}$. Let $l_{1} \in L$ such that $\psi\left(l_{1}\right)=m_{1}$. Then $l$ has to be replaced by an element $l^{\prime} \in L$ such that

$$
\psi\left(l^{\prime}\right)=m+d m_{1}=\psi\left(l+d l_{1}\right)
$$

It follows that $l^{\prime}=l+d l_{1}+\varphi\left(k_{1}\right)$ for some $k_{1} \in K$, hence

$$
d l^{\prime}=d l+d \varphi\left(k_{1}\right)=\varphi(k)+\varphi\left(d k_{1}\right)=\varphi\left(k^{\prime}\right),
$$

therefore $k^{\prime}=k+d k_{1}$ and $k^{\prime}$ has the same cohomology class as $k$.
Now, let us show that ker $\partial=\operatorname{Im} H(\psi)$. If $\{m\}$ is in the image of $H(\psi)$, we can take $m=\psi(l)$ with $d l=0$, thus $\partial\{m\}=0$. Conversely, if $\partial\{m\}=\{k\}=0$, we have $k=d k_{1}$ for some $k_{1} \in K$, hence $d l=\varphi(k)=d \varphi\left(k_{1}\right), z:=l-\varphi\left(k_{1}\right) \in Z(L)$ and $m=\psi(l)=\psi(z)$ is in $\operatorname{Im} H(\psi)$. We leave the verification of the other equalities $\operatorname{Im} H(\varphi)=\operatorname{ker} H(\psi)$, $\operatorname{Im} \partial=\operatorname{ker} H(\varphi)$ and of the commutation statement to the reader.

In most applications, the differential modules come with a natural $\mathbb{Z}$-grading. A homological complex is a graded differential module $K_{\bullet}=\bigoplus_{q \in \mathbb{Z}} K_{q}$ together with a differential $d$ of degree -1 , i.e. $d=\bigoplus d_{q}$ with $d_{q}: K_{q} \longrightarrow K_{q-1}$ and $d_{q-1} \circ d_{q}=$ 0 . Similarly, a cohomological complex is a graded differential module $K^{\bullet}=\bigoplus_{q \in \mathbb{Z}} K^{q}$ with differentials $d^{q}: K^{q} \longrightarrow K^{q+1}$ such that $d^{q+1} \circ d^{q}=0$ (superscripts are always used instead of subscripts in that case). The corresponding (co)cycle, (co)boundary and (co)homology modules inherit a natural $\mathbb{Z}$-grading. In the case of cohomology, say, these modules will be denoted

$$
Z^{\bullet}\left(K^{\bullet}\right)=\bigoplus Z^{q}\left(K^{\bullet}\right), \quad B^{\bullet}\left(K^{\bullet}\right)=\bigoplus B^{q}\left(K^{\bullet}\right), \quad H^{\bullet}\left(K^{\bullet}\right)=\bigoplus H^{q}\left(K^{\bullet}\right)
$$

Unless otherwise stated, morphisms of complexes are assumed to be of degree 0, i.e. of the form $\varphi^{\bullet}=\bigoplus \varphi^{q}$ with $\varphi^{q}: K^{q} \longrightarrow L^{q}$. Any short exact sequence

$$
0 \longrightarrow K^{\bullet} \xrightarrow{\varphi^{\bullet}} L^{\bullet} \xrightarrow{\psi^{\bullet}} M^{\bullet} \longrightarrow 0
$$

gives rise to a corresponding long exact sequence of cohomology groups

$$
\begin{equation*}
\left.H^{q}\left(K^{\bullet}\right) \xrightarrow{H^{q}\left(\varphi^{\bullet}\right)} H^{q}\left(L^{\bullet}\right) \xrightarrow{H^{q}\left(\psi^{\bullet}\right)} H^{q}\left(M^{\bullet}\right) \xrightarrow{\partial^{q}} H^{q+1}\left(K^{\bullet}\right) \xrightarrow{H^{q+1}(\varphi}\right) . \tag{1.6}
\end{equation*}
$$

and there is a similar homology long exact sequence with a connecting homomorphism $\partial_{q}$ of degree -1 . When dealing with commutative diagrams of such sequences, the following simple lemma is often useful; the proof consists in a straightforward diagram chasing.
(1.7) Five lemma. Consider a commutative diagram of $R$-modules

$$
\begin{array}{llll}
A_{1} \longrightarrow & A_{2} \longrightarrow A_{3} \longrightarrow A_{4} \longrightarrow A_{5} \\
\downarrow \varphi_{1} & \downarrow \varphi_{2} & \downarrow \varphi_{3} & \downarrow \varphi_{4} \\
B_{1} \longrightarrow & \downarrow \varphi_{5} \\
B_{2} \longrightarrow & B_{3} \longrightarrow & B_{4} \longrightarrow & B_{5}
\end{array}
$$

where the rows are exact sequences. If $\varphi_{2}$ and $\varphi_{4}$ are injective and $\varphi_{1}$ surjective, then $\varphi_{3}$ is injective. If $\varphi_{2}$ and $\varphi_{4}$ is surjective and $\varphi_{5}$ injective, then $\varphi_{3}$ is surjective. In particular, $\varphi_{3}$ is an isomorphism as soon as $\varphi_{1}, \varphi_{2}, \varphi_{4}, \varphi_{5}$ are isomorphisms.

## § 2. The Simplicial Flabby Resolution of a Sheaf

Let $X$ be a topological space and let $\mathcal{A}$ be a sheaf of abelian groups on $X$ (see $\S$ II-2 for the definition). All the sheaves appearing in the sequel are assumed implicitly to be sheaves of abelian groups, unless otherwise stated. The first useful notion is that of resolution.
(2.1) Definition. $A$ (cohomological) resolution of $\mathscr{A}$ is a differential complex of sheaves $\left(\mathscr{L}^{\bullet}, d\right)$ with $\mathscr{L}^{q}=0, d^{q}=0$ for $q<0$, such that there is an exact sequence

$$
0 \longrightarrow \mathscr{A} \xrightarrow{j} \mathscr{L}^{0} \xrightarrow{d^{0}} \mathscr{L}^{1} \longrightarrow \cdots \longrightarrow \mathscr{L}^{q} \xrightarrow{d^{q}} \mathscr{L}^{q+1} \longrightarrow \cdots .
$$

If $\varphi: \mathscr{A} \longrightarrow \mathscr{B}$ is a morphism of sheaves and $\left(\mathscr{M}^{\bullet}, d\right)$ a resolution of $\mathscr{B}$, a morphism of resolutions $\varphi^{\bullet}: \mathscr{L}^{\bullet} \longrightarrow \mathcal{M}^{\bullet}$ is a commutative diagram
(2.2) Example. Let $X$ be a differentiable manifold and $\mathscr{C}^{q}$ the sheaf of germs of $\mathscr{C}^{\infty}$ differential forms of degree $q$ with real values. The exterior derivative $d$ defines a resolution $\left(\mathscr{E}^{\bullet}, d\right)$ of the sheaf $\mathbb{R}$ of locally constant functions with real values. In fact Poincaré's lemma asserts that $d$ is locally exact in degree $q \geqslant 1$, and it is clear that the sections of $\operatorname{ker} d^{0}$ on connected open sets are constants.

In the sequel, we will be interested by special resolutions in which the sheaves $\mathscr{L}^{q}$ have no local "rigidity". For that purpose, we introduce flabby sheaves, which have become a standard tool in sheaf theory since the publication of Godement's book [Godement 1957].
(2.3) Definition. A sheaf $\mathscr{F}$ is called flabby if for every open subset $U$ of $X$, the restriction map $\mathscr{F}(X) \longrightarrow \mathscr{F}(U)$ is onto, i.e. if every section of $\mathscr{F}$ on $U$ can be extended to $X$.

Let $\pi: \mathscr{A} \longrightarrow X$ be a sheaf on $X$. We denote by $\mathscr{A}^{[0]}$ the sheaf of germs of sections $X \longrightarrow \mathscr{A}$ which are not necessarily continuous. In other words, $\mathscr{A}^{[0]}(U)$ is the set of all maps $f: U \longrightarrow \mathscr{A}$ such that $f(x) \in \mathscr{A}_{x}$ for all $x \in U$, or equivalently $\mathscr{A}^{[0]}(U)=\prod_{x \in U} \mathscr{A}_{x}$. It is clear that $\mathscr{A}^{[0]}$ is flabby and there is a canonical injection

$$
j: \mathscr{A} \longrightarrow \mathscr{A}^{[0]}
$$

defined as follows: to any $s \in \mathscr{A}_{x}$ we associate the germ $\widetilde{s} \in \mathscr{A}_{x}^{[0]}$ equal to the continuous section $y \longmapsto \widetilde{s}(y)$ near $x$ such that $\widetilde{s}(x)=s$. In the sequel we merely denote $\widetilde{s}: y \longmapsto s(y)$ for simplicity. The sheaf $\mathscr{A}^{[0]}$ is called the canonical flabby sheaf associated to $\mathscr{A}$. We define inductively

$$
\mathscr{A}^{[q]}=\left(\mathscr{A}^{[q-1]}\right)^{[0]} .
$$

The stalk $\mathscr{A}_{x}^{[q]}$ can be considered as the set of equivalence classes of maps $f: X^{q+1} \longrightarrow \mathscr{A}$ such that $f\left(x_{0}, \ldots, x_{q}\right) \in \mathscr{A}_{x_{q}}$, with two such maps identified if they coincide on a set of the form

$$
\begin{equation*}
x_{0} \in V, \quad x_{1} \in V\left(x_{0}\right), \quad \ldots, \quad x_{q} \in V\left(x_{0}, \ldots, x_{q-1}\right), \tag{2.4}
\end{equation*}
$$

where $V$ is an open neighborhood of $x$ and $V\left(x_{0}, \ldots, x_{j}\right)$ an open neighborhood of $x_{j}$, depending on $x_{0}, \ldots, x_{j}$. This is easily seen by induction on $q$, if we identify a $\operatorname{map} f: X^{q+1} \rightarrow \mathscr{A}$ to the map $X \rightarrow \mathscr{A}^{[q-1]}, x_{0} \mapsto f_{x_{0}}$ such that $f_{x_{0}}\left(x_{1}, \ldots, x_{q}\right)=$ $f\left(x_{0}, x_{1}, \ldots, x_{q}\right)$. Similarly, $\mathscr{A}^{[q]}(U)$ is the set of equivalence classes of functions $X^{q+1} \ni$ $\left(x_{0}, \ldots, x_{q}\right) \longmapsto f\left(x_{0}, \ldots, x_{q}\right) \in \mathscr{A}_{x_{q}}$, with two such functions identified if they coincide on a set of the form

$$
x_{0} \in U, \quad x_{1} \in V\left(x_{0}\right), \quad \ldots, \quad x_{q} \in V\left(x_{0}, \ldots, x_{q-1}\right)
$$

Here, we may of course suppose $V\left(x_{0}, \ldots, x_{q-1}\right) \subset \ldots \subset V\left(x_{0}, x_{1}\right) \subset V\left(x_{0}\right) \subset U$. We define a differential $d^{q}: \mathscr{A}^{[q]} \longrightarrow \mathscr{A}^{[q+1]}$ by

$$
\begin{align*}
& \left(d^{q} f\right)\left(x_{0}, \ldots, x_{q+1}\right)=  \tag{2.5}\\
& \sum_{0 \leqslant j \leqslant q}(-1)^{j} f\left(x_{0}, \ldots, \widehat{x_{j}}, \ldots, x_{q+1}\right)+(-1)^{q+1} f\left(x_{0}, \ldots, x_{q}\right)\left(x_{q+1}\right) .
\end{align*}
$$

The meaning of the last term is to be understood as follows: the element $s=f\left(x_{0}, \ldots, x_{q}\right)$ is a germ in $\mathscr{A}_{x_{q}}$, therefore $s$ defines a continuous section $x_{q+1} \mapsto s\left(x_{q+1}\right)$ of $\mathscr{A}$ in a neighborhood $V\left(x_{0}, \ldots, x_{q}\right)$ of $x_{q}$. In low degrees, we have the formulas

$$
\begin{align*}
(j s)\left(x_{0}\right) & =s\left(x_{0}\right), \quad s \in \mathscr{A}_{x}, \\
\left(d^{0} f\right)\left(x_{0}, x_{1}\right) & =f\left(x_{1}\right)-f\left(x_{0}\right)\left(x_{1}\right), \quad f \in \mathscr{A}_{x}^{[0]},  \tag{2.6}\\
\left(d^{1} f\right)\left(x_{0}, x_{1}, x_{2}\right) & =f\left(x_{1}, x_{2}\right)-f\left(x_{0}, x_{2}\right)+f\left(x_{0}, x_{1}\right)\left(x_{2}\right), \quad f \in \mathscr{A}_{x}^{[1]} .
\end{align*}
$$

(2.7) Theorem ([Godement 1957]). The complex $\left(\mathscr{A}^{[\bullet]}, d\right)$ is a resolution of the sheaf $\mathscr{A}$, called the simplicial flabby resolution of $\mathscr{A}$.

Proof. For $s \in \mathscr{A}_{x}$, the associated continuous germ obviously satisfies $s\left(x_{0}\right)\left(x_{1}\right)=s\left(x_{1}\right)$ for $x_{0} \in V, x_{1} \in V\left(x_{0}\right)$ small enough. The reader will easily infer from this that $d^{0} \circ j=0$ and $d^{q+1} \circ d^{q}=0$. In order to verify that $\left(\mathscr{A}^{[\bullet]}, d\right)$ is a resolution of $\mathscr{A}$, we show that the complex

$$
\cdots \longrightarrow 0 \longrightarrow \mathscr{A}_{x} \xrightarrow{j} \mathscr{A}_{x}^{[0]} \xrightarrow{d^{0}} \cdots \longrightarrow \mathscr{A}_{x}^{[q]} \xrightarrow{d^{q}} \mathscr{A}_{x}^{[q+1]} \longrightarrow \cdots
$$

is homotopic to zero for every point $x \in X$. Set $\mathscr{A}^{[-1]}=\mathscr{A}, d^{-1}=j$ and

$$
\begin{array}{lll}
h^{0}: & \mathscr{A}_{x}^{[0]} \longrightarrow \mathscr{A}_{x}, & h^{0}(f)=f(x) \in \mathscr{A}_{x}, \\
h^{q}: & \mathscr{A}_{x}^{[q]} \longrightarrow \mathscr{A}_{x}^{[q-1]}, & h^{q}(f)\left(x_{0}, \ldots, x_{q-1}\right)=f\left(x, x_{0}, \ldots, x_{q-1}\right) .
\end{array}
$$

A straightforward computation shows that $\left(h^{q+1} \circ d^{q}+d^{q-1} \circ h^{q}\right)(f)=f$ for all $q \in \mathbb{Z}$ and $f \in A_{x}^{[q]}$.

If $\varphi: \mathscr{A} \longrightarrow \mathscr{B}$ is a sheaf morphism, it is clear that $\varphi$ induces a morphism of resolutions

$$
\begin{equation*}
\varphi^{[\bullet]}: \mathscr{A}[\bullet] \longrightarrow \mathscr{B}^{[\bullet]} \tag{2.8}
\end{equation*}
$$

For every short exact sequence $\mathscr{A} \rightarrow \mathscr{B} \rightarrow \mathscr{C}$ of sheaves, we get a corresponding short exact sequence of sheaf complexes

$$
\begin{equation*}
\mathscr{A}^{[\bullet]} \longrightarrow \mathscr{B}^{[\bullet]} \longrightarrow \mathscr{C}^{[\bullet]} . \tag{2.9}
\end{equation*}
$$

## § 3. Cohomology Groups with Values in a Sheaf

## § 3.A. Definition and Functorial Properties

If $\pi: \mathscr{A} \rightarrow X$ is a sheaf of abelian groups, the cohomology groups of $\mathscr{A}$ on $X$ are (in a vague sense) algebraic invariants which describe the rigidity properties of the global sections of $\mathscr{A}$.
(3.1) Definition. For every $q \in \mathbb{Z}$, the $q$-th cohomology group of $X$ with values in $\mathscr{A}$ is

$$
\begin{aligned}
& H^{q}(X, \mathscr{A})=H^{q}\left(\mathscr{A}^{[\bullet]}(X)\right)= \\
& \quad=\operatorname{ker}\left(d^{q}: \mathscr{A}^{[q]}(X) \rightarrow \mathscr{A}^{[q+1]}(X)\right) / \operatorname{Im}\left(d^{q-1}: \mathscr{A}^{[q-1]}(X) \rightarrow \mathscr{A}^{[q]}(X)\right)
\end{aligned}
$$

with the convention $\mathscr{A}^{[q]}=0, d^{q}=0, H^{q}(X, \mathscr{A})=0$ when $q<0$.

For any subset $S \subset X$, we denote by $\mathcal{A}_{\mid S}$ the restriction of $\mathscr{A}$ to $S$, i.e. the sheaf $\mathscr{A}_{\mid S}=\pi^{-1}(S)$ equipped with the projection $\pi_{\upharpoonright S}$ onto $S$. Then we write $H^{q}\left(S, \mathscr{A}_{\uparrow S}\right)=$ $H^{q}(S, \mathscr{A})$ for simplicity. When $U$ is open, we see that $\left(\mathscr{A}^{[q]}\right)_{\mid U}$ coincides with $\left(\mathscr{A}_{\mid U}\right)^{[q]}$, thus we have $H^{q}(U, \mathscr{A})=H^{q}\left(\mathscr{A}^{[\bullet]}(U)\right)$. It is easy to show that every exact sequence of sheaves $0 \rightarrow \mathscr{A} \rightarrow \mathscr{L}^{0} \rightarrow \mathscr{L}^{1}$ induces an exact sequence

$$
\begin{equation*}
0 \longrightarrow \mathscr{A}(X) \longrightarrow \mathscr{L}^{0}(X) \longrightarrow \mathscr{L}^{1}(X) \tag{3.2}
\end{equation*}
$$

If we apply this to $\mathscr{L}^{q}=\mathscr{A}^{[q]}, q=0,1$, we conclude that

$$
\begin{equation*}
H^{0}(X, \mathscr{A})=\mathscr{A}(X) \tag{3.3}
\end{equation*}
$$

Let $\varphi: \mathscr{A} \longrightarrow \mathscr{B}$ be a sheaf morphism; (2.8) shows that there is an induced morphism

$$
\begin{equation*}
H^{q}(\varphi): H^{q}(X, \mathscr{A}) \longrightarrow H^{q}(X, \mathscr{B}) \tag{3.4}
\end{equation*}
$$

on cohomology groups. Let $0 \rightarrow \mathscr{A} \rightarrow \mathscr{B} \rightarrow \mathscr{C} \rightarrow 0$ be an exact sequence of sheaves. Then we have an exact sequence of groups

$$
0 \longrightarrow \mathscr{A}^{[0]}(X) \longrightarrow \mathscr{B}^{[0]}(X) \longrightarrow \mathscr{C}^{[0]}(X) \longrightarrow 0
$$

because $\mathscr{A}^{[0]}(X)=\prod_{x \in X} \mathscr{A}_{x}$. Similarly, (2.9) yields for every $q$ an exact sequence of groups

$$
0 \longrightarrow \mathscr{A}^{[q]}(X) \longrightarrow \mathscr{B}^{[q]}(X) \longrightarrow \mathscr{C}^{[q]}(X) \longrightarrow 0
$$

If we take (3.3) into account, the snake lemma implies:
(3.5) Theorem. To any exact sequence of sheaves $0 \rightarrow \mathscr{A} \rightarrow \mathscr{B} \rightarrow \mathscr{C} \rightarrow 0$ is associated a long exact sequence of cohomology groups

$$
\begin{aligned}
& 0 \longrightarrow \mathscr{A}(X) \longrightarrow \mathscr{B}(X) \\
& \cdots \longrightarrow \mathscr{C}(X) \longrightarrow H^{1}(X, \mathscr{A}) \longrightarrow \cdots \\
& \cdots H^{q}(X, \mathscr{A}) \longrightarrow H^{q}(X, \mathscr{B}) \longrightarrow H^{q}(X, \mathscr{C}) \longrightarrow H^{q+1}(X, \mathscr{A}) \longrightarrow \cdots .
\end{aligned}
$$

(3.6) Corollary. Let $\mathscr{B} \rightarrow \mathscr{C}$ be a surjective sheaf morphism and let $\mathscr{A}$ be its kernel. If $H^{1}(X, \mathscr{A})=0$, then $\mathscr{B}(X) \rightarrow \mathscr{C}(X)$ is surjective.

## § 3.B. Exact Sequence Associated to a Closed Subset

Let $S$ be a closed subset of $X$ and $U=X \backslash S$. For any sheaf $\mathscr{A}$ on $X$, the presheaf

$$
\Omega \longmapsto \mathscr{A}(S \cap \Omega), \quad \Omega \subset X \text { open }
$$

with the obvious restriction maps satisfies axioms (II-2.4') and (II-2.4"), so it defines a sheaf on $X$ which we denote by $\mathscr{A}^{S}$. This sheaf should not be confused with the restriction sheaf $\mathscr{A}_{\Gamma S}$, which is a sheaf on $S$. We easily find

$$
\begin{equation*}
\left(\mathscr{A}^{S}\right)_{x}=\mathscr{A}_{x} \quad \text { if } \quad x \in S, \quad\left(\mathscr{A}^{S}\right)_{x}=0 \quad \text { if } \quad x \in U \tag{3.7}
\end{equation*}
$$

Observe that these relations would completely fail if $S$ were not closed. The restriction morphism $f \mapsto f_{\lceil S}$ induces a surjective sheaf morphism $\mathscr{A} \rightarrow \mathscr{A}^{S}$. We let $\mathscr{A}_{U}$ be its kernel, so that we have the relations

$$
\begin{equation*}
\left(\mathscr{A}_{U}\right)_{x}=0 \quad \text { if } x \in S, \quad\left(\mathscr{A}_{U}\right)_{x}=\mathscr{A}_{x} \quad \text { if } \quad x \in U . \tag{3.8}
\end{equation*}
$$

From the definition, we obtain in particular

$$
\begin{equation*}
\mathscr{A}^{S}(X)=\mathscr{A}(S), \quad \mathscr{A}_{U}(X)=\{\text { sections of } \mathscr{A}(X) \text { vanishing on } S\} . \tag{3.9}
\end{equation*}
$$

Theorem 3.5 applied to the exact sequence $0 \rightarrow \mathscr{A}_{U} \rightarrow \mathscr{A} \rightarrow \mathscr{A}^{S} \rightarrow 0$ on $X$ gives a long exact sequence

$$
\begin{align*}
0 & \longrightarrow A_{U}(X) \tag{3.9}
\end{align*} \longrightarrow \mathscr{A}(X) \longrightarrow \mathscr{A}(S) \longrightarrow H^{1}\left(X, \mathscr{A}_{U}\right) \cdots .
$$

## § 3.C. Mayer-Vietoris Exact Sequence

Let $U_{1}, U_{2}$ be open subsets of $X$ and $U=U_{1} \cup U_{2}, V=U_{1} \cap U_{2}$. For any sheaf $\mathscr{A}$ on $X$ and any $q$ we have an exact sequence

$$
0 \longrightarrow \mathscr{A}^{[q]}(U) \longrightarrow \mathscr{A}^{[q]}\left(U_{1}\right) \oplus \mathscr{A}^{[q]}\left(U_{2}\right) \longrightarrow \mathscr{A}^{[q]}(V) \longrightarrow 0
$$

where the injection is given by $f \longmapsto\left(f_{\mid U_{1}}, f_{\mid U_{2}}\right)$ and the surjection by $\left(g_{1}, g_{2}\right) \longmapsto$ $g_{2 \upharpoonright V}-g_{1 \upharpoonright V}$; the surjectivity of this map follows immediately from the fact that $\mathcal{A}^{[q]}$ is flabby. An application of the snake lemma yields:
(3.11) Theorem. For any sheaf $\mathscr{A}$ on $X$ and any open sets $U_{1}, U_{2} \subset X$, set $U=U_{1} \cup U_{2}$, $V=U_{1} \cap U_{2}$. Then there is an exact sequence

$$
H^{q}(U, \mathscr{A}) \longrightarrow H^{q}\left(U_{1}, \mathscr{A}\right) \oplus H^{q}\left(U_{2}, \mathscr{A}\right) \longrightarrow H^{q}(V, \mathscr{A}) \longrightarrow H^{q+1}(U, \mathscr{A}) \cdots
$$

## §4. Acyclic Sheaves

Given a sheaf $\mathscr{A}$ on $X$, it is usually very important to decide whether the cohomology groups $H^{q}(U, \mathcal{A})$ vanish for $q \geqslant 1$, and if this is the case, for which type of open sets $U$. Note that one cannot expect to have $H^{0}(U, \mathscr{A})=0$ in general, since a sheaf always has local sections.
(4.1) Definition. $A$ sheaf $\mathscr{A}$ is said to be acyclic on an open subset $U$ if $H^{q}(U, \mathscr{A})=0$ for $q \geqslant 1$.

## $\S$ 4.A. Case of Flabby Sheaves

We are going to show that flabby sheaves are acyclic. First we need the following simple result.
(4.2) Proposition. Let $\mathscr{A}$ be a sheaf with the following property: for every section $f$ of $\mathscr{A}$ on an open subset $U \subset X$ and every point $x \in X$, there exists a neighborhood $\Omega$ of $x$ and a section $h \in \mathscr{A}(\Omega)$ such that $h=f$ on $U \cap \Omega$. Then $\mathscr{A}$ is flabby.

A consequence of this proposition is that flabbiness is a local property: a sheaf $\mathscr{A}$ is flabby on $X$ if and only if it is flabby on a neighborhood of every point of $X$.

Proof. Let $f \in \mathscr{A}(U)$ be given. Consider the set of pairs $(v, V)$ where $v$ in $\mathscr{B}(V)$ is an extension of $f$ on an open subset $V \supset U$. This set is inductively ordered, so there exists a maximal extension $(v, V)$ by Zorn's lemma. The assumption shows that $V$ must be equal to $X$.
(4.3) Proposition. Let $0 \longrightarrow \mathscr{A} \xrightarrow{j} \mathscr{B} \xrightarrow{p} \mathscr{C} \longrightarrow 0$ be an exact sequence of sheaves. If $\mathscr{A}$ is flabby, the sequence of groups

$$
0 \longrightarrow \mathscr{A}(U) \xrightarrow{j} \mathscr{B}(U) \xrightarrow{p} \mathscr{C}(U) \longrightarrow 0
$$

is exact for every open set $U$. If $\mathscr{A}$ and $\mathscr{B}$ are flabby, then $\mathscr{C}$ is flabby.
Proof. Let $g \in \mathscr{C}(U)$ be given. Consider the set $E$ of pairs $(v, V)$ where $V$ is an open subset of $U$ and $v \in \mathscr{B}(V)$ is such that $p(v)=g$ on $V$. It is clear that $E$ is inductively ordered, so $E$ has a maximal element $(v, V)$, and we will prove that $V=U$. Otherwise, let $x \in U \backslash V$ and let $h$ be a section of $\mathscr{B}$ in a neighborhood of $x$ such that $p\left(h_{x}\right)=g_{x}$. Then $p(h)=g$ on a neighborhood $\Omega$ of $x$, thus $p(v-h)=0$ on $V \cap \Omega$ and $v-h=j(u)$ with $u \in \mathscr{A}(V \cap \Omega)$. If $\mathscr{A}$ is flabby, $u$ has an extension $\widetilde{u} \in \mathscr{A}(X)$ and we can define a section $w \in \mathscr{B}(V \cup \Omega)$ such that $p(w)=g$ by

$$
w=v \quad \text { on } V, \quad w=h+j(\widetilde{u}) \text { on } \Omega,
$$

contradicting the maximality of $(v, V)$. Therefore $V=U, v \in \mathscr{B}(U)$ and $p(v)=g$ on $U$. The first statement is proved. If $\mathscr{B}$ is also flabby, $v$ has an extension $\widetilde{v} \in \mathscr{B}(X)$ and $\widetilde{g}=p(\widetilde{v}) \in \mathscr{C}(X)$ is an extension of $g$. Hence $\mathscr{C}$ is flabby.
(4.4) Theorem. $A$ flabby sheaf $\mathscr{A}$ is acyclic on all open sets $U \subset X$.

Proof. Let $\mathscr{\mathscr { L } ^ { q }}=\operatorname{ker}\left(d^{q}: \mathscr{A}^{[q]} \rightarrow \mathscr{A}^{[q+1]}\right)$. Then $\mathscr{L}^{0}=\mathscr{A}$ and we have an exact sequence of sheaves

$$
0 \longrightarrow \mathscr{Z}^{q} \longrightarrow \mathscr{A}^{[q]} \xrightarrow{d^{q}} \mathscr{Z}^{q+1} \longrightarrow 0
$$

because $\operatorname{Im} d^{q}=\operatorname{ker} d^{q+1}=\mathscr{L}^{q+1}$. Proposition 4.3 implies by induction on $q$ that all sheaves $\mathscr{L}^{q}$ are flabby, and yields exact sequences

$$
0 \longrightarrow \mathscr{Z}^{q}(U) \longrightarrow \mathscr{A}^{[q]}(U) \xrightarrow{d^{q}} \mathscr{Z}^{q+1}(U) \longrightarrow 0 .
$$

For $q \geqslant 1$, we find therefore

$$
\begin{aligned}
\operatorname{ker}\left(d^{q}: \mathscr{A}^{[q]}(U) \rightarrow \mathscr{A}^{[q+1]}(U)\right) & =\mathscr{\not} \mathscr{L}^{q}(U) \\
& =\operatorname{Im}\left(d^{q-1}: \mathscr{A}^{[q-1]}(U) \rightarrow \mathscr{A}^{[q]}(U)\right),
\end{aligned}
$$

that is, $H^{q}(U, \mathscr{A})=H^{q}\left(\mathscr{A}^{[\bullet]}(U)\right)=0$.

## §4.B. Soft Sheaves over Paracompact Spaces

We now discuss another general situation which produces acyclic sheaves. Recall that a topological space $X$ is said to be paracompact if $X$ is Hausdorff and if every open covering of $X$ has a locally finite refinement. For instance, it is well known that every metric space is paracompact. A paracompact space $X$ is always normal; in particular, for any locally finite open covering $\left(U_{\alpha}\right)$ of $X$ there exists an open covering $\left(V_{\alpha}\right)$ such that $\bar{V}_{\alpha} \subset U_{\alpha}$. We will also need another closely related concept.
(4.5) Definition. We say that a subspace $S$ is strongly paracompact in $X$ if $S$ is Hausdorff and if the following property is satisfied: for every covering $\left(U_{\alpha}\right)$ of $S$ by open sets in $X$, there exists another such covering $\left(V_{\beta}\right)$ and a neighborhood $W$ of $S$ such that each set $W \cap \bar{V}_{\beta}$ is contained in some $U_{\alpha}$, and such that every point of $S$ has a neighborhood intersecting only finitely many sets $V_{\beta}$.

It is clear that a strongly paracompact subspace $S$ is itself paracompact. Conversely, the following result is easy to check:
(4.6) Lemma. A subspace $S$ is strongly paracompact in $X$ as soon as one of the following situations occurs:
a) $X$ is paracompact and $S$ is closed;
b) $S$ has a fundamental family of paracompact neighborhoods in $X$;
c) $S$ is paracompact and has a neighborhood homeomorphic to some product $S \times T$, in which $S$ is embedded as a slice $S \times\left\{t_{0}\right\}$.
(4.7) Theorem. Let $\mathscr{A}$ be a sheaf on $X$ and $S$ a strongly paracompact subspace of $X$. Then every section $f$ of $\mathscr{A}$ on $S$ can be extended to a section of $\mathscr{A}$ on some open neighborhood $\Omega$ of $\mathcal{A}$.

Proof. Let $f \in \mathscr{A}(S)$. For every point $z \in S$ there exists an open neighborhood $U_{z}$ and a section $\widetilde{f}_{z} \in \mathscr{A}\left(U_{z}\right)$ such that $\widetilde{f}_{z}(z)=f(z)$. After shrinking $U_{z}$, we may assume that $\widetilde{f}_{z}$
and $f$ coincide on $S \cap U_{z}$. Let $\left(V_{\alpha}\right)$ be an open covering of $S$ that is locally finite near $S$ and $W$ a neighborhood of $S$ such that $W \cap \bar{V}_{\alpha} \subset U_{z(\alpha)}$ (Def. 4.5). We let

$$
\Omega=\left\{x \in W \cap \bigcup V_{\alpha} ; \tilde{f}_{z(\alpha)}(x)=\tilde{f}_{z(\beta)}(x), \forall \alpha, \beta \text { with } x \in \bar{V}_{\alpha} \cap \bar{V}_{\beta}\right\} .
$$

Then $\left(\Omega \cap V_{\alpha}\right)$ is an open covering of $\Omega$ and all pairs of sections $\tilde{f}_{z(\alpha)}$ coincide in pairwise intersections. Thus there exists a section $F$ of $\mathscr{A}$ on $\Omega$ which is equal to $\widetilde{f}_{z(\alpha)}$ on $\Omega \cap V_{\alpha}$. It remains only to show that $\Omega$ is a neighborhood of $S$. Let $z_{0} \in S$. There exists a neighborhood $U^{\prime}$ of $z_{0}$ which meets only finitely many sets $V_{\alpha_{1}}, \ldots, V_{\alpha_{p}}$. After shrinking $U^{\prime}$, we may keep only those $V_{\alpha_{l}}$ such that $z_{0} \in \bar{V}_{\alpha_{l}}$. The sections $\widetilde{f}_{z\left(\alpha_{l}\right)}$ coincide at $z_{0}$, so they coincide on some neighborhood $U^{\prime \prime}$ of this point. Hence $W \cap U^{\prime \prime} \subset \Omega$, so $\Omega$ is a neighborhood of $S$.
(4.8) Corollary. If $X$ is paracompact, every section $f \in \mathscr{A}(S)$ defined on a closed set $S$ extends to a neighborhood $\Omega$ of $S$.
(4.9) Definition. A sheaf $\mathscr{A}$ on $X$ is said to be soft if every section $f$ of $\mathscr{A}$ on a closed set $S$ can be extended to $X$, i.e. if the restriction map $\mathcal{A}(X) \longrightarrow \mathscr{A}(S)$ is onto for every closed set $S$.
(4.10) Example. On a paracompact space, every flabby sheaf $\mathscr{A}$ is soft: this is a consequence of Cor. 4.8.
(4.11) Example. On a paracompact space, the Tietze-Urysohn extension theorem shows that the sheaf $\mathscr{C}_{X}$ of germs of continuous functions on $X$ is a soft sheaf of rings. However, observe that $\mathscr{C}_{X}$ is not flabby as soon as $X$ is not discrete.
(4.12) Example. If $X$ is a paracompact differentiable manifold, the sheaf $\mathscr{E}_{X}$ of germs of $\mathscr{C}^{\infty}$ functions on $X$ is a soft sheaf of rings.

Until the end of this section, we assume that $X$ is a paracompact topological space. We first show that softness is a local property.
(4.13) Proposition. A sheaf $\mathscr{A}$ is soft on $X$ if and only if it is soft in a neighborhood of every point $x \in X$.

Proof. If $\mathscr{A}$ is soft on $X$, it is soft on any closed neighborhood of a given point. Conversely, let $\left(U_{\alpha}\right)_{\alpha \in I}$ be a locally finite open covering of $X$ which refine some covering by neighborhoods on which $\mathscr{A}$ is soft. Let $\left(V_{\alpha}\right)$ be a finer covering such that $\bar{V}_{\alpha} \subset U_{\alpha}$, and $f \in \mathscr{A}(S)$ be a section of $\mathscr{A}$ on a closed subset $S$ of $X$. We consider the set $E$ of pairs $(g, J)$, where $J \subset I$ and where $g$ is a section over $F_{J}:=S \cup \bigcup_{\alpha \in J} \bar{V}_{\alpha}$, such that $g=f$ on $S$. As the family $\left(\bar{V}_{\alpha}\right)$ is locally finite, a section of $\mathscr{A}$ over $F_{J}$ is continuous as soon it is continuous on $S$ and on each $\bar{V}_{\alpha}$. Then $(f, \emptyset) \in E$ and $E$ is inductively ordered by the relation

$$
\left(g^{\prime}, J^{\prime}\right) \longrightarrow\left(g^{\prime \prime}, J^{\prime \prime}\right) \text { if } J^{\prime} \subset J^{\prime \prime} \text { and } g^{\prime}=g^{\prime \prime} \text { on } F_{J^{\prime}}
$$

No element $(\underline{g}, J), J \neq I$, can be maximal: the assumption shows that $g_{\mid F_{J} \cap \bar{V}_{\alpha}}$ has an extension to $\bar{V}_{\alpha}$, thus such a $g$ has an extension to $F_{J \cup\{\alpha\}}$ for any $\alpha \notin J$. Hence $E$ has a maximal element $(g, I)$ defined on $F_{I}=X$.
(4.14) Proposition. Let $0 \rightarrow \mathscr{A} \rightarrow \mathscr{B} \rightarrow \mathscr{C} \rightarrow 0$ be an exact sequence of sheaves. If $\mathscr{A}$ is soft, the map $\mathscr{B}(S) \rightarrow \mathscr{C}(S)$ is onto for any closed subset $S$ of $X$. If $\mathscr{A}$ and $\mathscr{B}$ are soft, then $\mathscr{C}$ is soft.

By the above inductive method, this result can be proved in a way similar to its analogue for flabby sheaves. We therefore obtain:
(4.15) Theorem. On a paracompact space, a soft sheaf is acyclic on all closed subsets.
(4.16) Definition. The support of a section $f \in \mathscr{A}(X)$ is defined by

$$
\text { Supp } f=\{x \in X ; f(x) \neq 0\} .
$$

Supp $f$ is always a closed set: as $\mathscr{A} \rightarrow X$ is a local homeomorphism, the equality $f(x)=0$ implies $f=0$ in a neighborhood of $x$.
(4.17) Theorem. Let $\left(U_{\alpha}\right)_{\alpha \in I}$ be an open covering of $X$. If $\mathscr{A}$ is soft and $f \in \mathscr{A}(X)$, there exists a partition of $f$ subordinate to $\left(U_{\alpha}\right)$, i.e. a family of sections $f_{\alpha} \in \mathscr{A}(X)$ such that (Supp $f_{\alpha}$ ) is locally finite, Supp $f_{\alpha} \subset U_{\alpha}$ and $\sum f_{\alpha}=f$ on $X$.

Proof. Assume first that $\left(U_{\alpha}\right)$ is locally finite. There exists an open covering $\left(V_{\alpha}\right)$ such that $\bar{V}_{\alpha} \subset U_{\alpha}$. Let $\left(f_{\alpha}\right)_{\alpha \in J}, J \subset I$, be a maximal family of sections $f_{\alpha} \in \mathscr{A}(X)$ such that Supp $f_{\alpha} \subset U_{\alpha}$ and $\sum_{\alpha \in J} f_{\alpha}=f$ on $S=\bigcup_{\alpha \in J} \bar{V}_{\alpha}$. If $J \neq I$ and $\beta \in I \backslash J$, there exists a section $f_{\beta} \in \mathscr{A}(X)$ such that

$$
f_{\beta}=0 \quad \text { on } \quad X \backslash U_{\beta} \quad \text { and } \quad f_{\beta}=f-\sum_{\alpha \in J} f_{\alpha} \quad \text { on } \quad S \cup \bar{V}_{\beta}
$$

because $\left(X \backslash U_{\beta}\right) \cup S \cup \bar{V}_{\beta}$ is closed and $f-\sum f_{\alpha}=0$ on $\left(X \backslash U_{\alpha}\right) \cap S$. This is a contradiction unless $J=I$.

In general, let $\left(V_{j}\right)$ be a locally finite refinement of $\left(U_{\alpha}\right)$, such that $V_{j} \subset U_{\rho(j)}$, and let $\left(f_{j}^{\prime}\right)$ be a partition of $f$ subordinate to $\left(V_{j}\right)$. Then $f_{\alpha}=\sum_{j \in \rho^{-1}(\alpha)} f_{j}^{\prime}$ is the required partition of $f$.

Finally, we discuss a special situation which occurs very often in practice. Let $\mathscr{R}$ be a sheaf of commutative rings on $X$; the rings $\mathscr{R}_{x}$ are supposed to have a unit element. Assume that $\mathscr{A}$ is a sheaf of modules over $\mathscr{R}$. It is clear that $\mathscr{A}^{[0]}$ is a $\mathscr{R}^{[0]}$-module, and thus also a $\mathscr{R}$-module. Therefore all sheaves $\mathscr{A}^{[q]}$ are $\mathscr{R}$-modules and the cohomology groups $H^{q}(U, \mathscr{A})$ have a natural structure of $\mathscr{R}(U)$-module.
(4.18) Lemma. If $\mathscr{R}$ is soft, every sheaf $\mathscr{A}$ of $\mathscr{R}$-modules is soft.

Proof. Every section $f \in \mathscr{A}(S)$ defined on a closed set $S$ has an extension to some open neighborhood $\Omega$. Let $\psi \in \mathscr{R}(X)$ be such that $\psi=1$ on $S$ and $\psi=0$ on $X \backslash \Omega$. Then $\psi f$, defined as 0 on $X \backslash \Omega$, is an extension of $f$ to $X$.
(4.19) Corollary. Let $\mathscr{A}$ be a sheaf of $\mathscr{E}_{X}$-modules on a paracompact differentiable manifold $X$. Then $H^{q}(X, \mathscr{A})=0$ for all $q \geqslant 1$.

## § 5. Čech Cohomology

## §5.A. Definitions

In many important circumstances, cohomology groups with values in a sheaf $\mathscr{A}$ can be computed by means of the complex of Čech cochains, which is directly related to the spaces of sections of $\mathscr{A}$ on sufficiently fine coverings of $X$. This more concrete approach was historically the first one used to define sheaf cohomology ([Leray 1950], [Cartan 1950]); however Čech cohomology does not always coincide with the "good" cohomology on non paracompact spaces. Let $\mathscr{U}=\left(U_{\alpha}\right)_{\alpha \in I}$ be an open covering of $X$. For the sake of simplicity, we denote

$$
U_{\alpha_{0} \alpha_{1} \ldots \alpha_{q}}=U_{\alpha_{0}} \cap U_{\alpha_{1}} \cap \ldots \cap U_{\alpha_{q}} .
$$

The group $C^{q}(\mathscr{U}, \mathscr{A})$ of $\check{C}$ ech $q$-cochains is the set of families

$$
c=\left(c_{\alpha_{0} \alpha_{1} \ldots \alpha_{q}}\right) \in \prod_{\left(\alpha_{0}, \ldots, \alpha_{q}\right) \in I^{q+1}} \mathscr{A}\left(U_{\alpha_{0} \alpha_{1} \ldots \alpha_{q}}\right) .
$$

The group structure on $C^{q}(\mathscr{U}, \mathscr{A})$ is the obvious one deduced from the addition law on sections of $\mathscr{A}$. The Čech differential $\delta^{q}: C^{q}(\mathscr{U}, \mathscr{A}) \longrightarrow C^{q+1}(\mathscr{U}, \mathscr{A})$ is defined by the formula

$$
\begin{equation*}
\left(\delta^{q} c\right)_{\alpha_{0} \ldots \alpha_{q+1}}=\sum_{0 \leqslant j \leqslant q+1}(-1)^{j} c_{\alpha_{0} \ldots \widehat{\alpha_{j} \ldots \alpha_{q+1}} \mid U_{\alpha_{0} \ldots \alpha_{q+1}}, ~, ~, ~}, \tag{5.1}
\end{equation*}
$$

and we set $C^{q}(\mathscr{U}, \mathscr{A})=0, \delta^{q}=0$ for $q<0$. In degrees 0 and 1 , we get for example

$$
\begin{array}{ll}
q=0, & c=\left(c_{\alpha}\right), \quad\left(\delta^{0} c\right)_{\alpha \beta}=c_{\beta}-c_{\alpha \mid U_{\alpha \beta}}, \\
q=1, & c=\left(c_{\alpha \beta}\right), \quad\left(\delta^{1} c\right)_{\alpha \beta \gamma}=c_{\beta \gamma}-c_{\alpha \gamma}+c_{\alpha \beta \mid U_{\alpha \beta \gamma}} .
\end{array}
$$

Easy verifications left to the reader show that $\delta^{q+1} \circ \delta^{q}=0$. We get therefore a cochain complex $(C \bullet(\mathscr{U}, \mathcal{A}), \delta)$, called the complex of Čech cochains relative to the covering $\mathscr{U}$.
(5.3) Definition. The Čech cohomology group of $\mathscr{A}$ relative to $\mathscr{U}$ is

$$
\check{H}^{q}(\mathscr{U}, \mathscr{A})=H^{q}\left(C^{\bullet}(\mathscr{U}, \mathscr{A})\right) .
$$

Formula (5.2) shows that the set of Čech 0-cocycles is the set of families $\left(c_{\alpha}\right) \in$ $\prod \mathscr{A}\left(U_{\alpha}\right)$ such that $c_{\beta}=c_{\alpha}$ on $U_{\alpha} \cap U_{\beta}$. Such a family defines in a unique way a global section $f \in \mathscr{A}(X)$ with $f_{\left\lceil U_{\alpha}\right.}=c_{\alpha}$. Hence

$$
\begin{equation*}
\check{H}^{0}(\ddots, \mathscr{A})=\mathscr{A}(X) . \tag{5.4}
\end{equation*}
$$

Now, let $\mathscr{V}=\left(V_{\beta}\right)_{\beta \in J}$ be another open covering of $X$ that is finer than $\mathscr{U}$; this means that there exists a map $\rho: J \rightarrow I$ such that $V_{\beta} \subset U_{\rho(\beta)}$ for every $\beta \in J$. Then we can define a morphism $\rho^{\bullet}: C^{\bullet}(\mathscr{U}, \mathscr{A}) \longrightarrow C^{\bullet}(\mathscr{V}, \mathscr{A})$ by

$$
\begin{equation*}
\left(\rho^{q} c\right)_{\beta_{0} \ldots \beta_{q}}=c_{\rho\left(\beta_{0}\right) \ldots \rho\left(\beta_{q}\right) \mid V_{\beta_{0} \ldots \beta_{q}}} ; \tag{5.5}
\end{equation*}
$$

the commutation property $\delta \rho^{\bullet}=\rho^{\bullet} \delta$ is immediate. If $\rho^{\prime}: J \rightarrow I$ is another refinement map such that $V_{\beta} \subset U_{\rho^{\prime}(\beta)}$ for all $\beta$, the morphisms $\rho^{\bullet}, \rho^{\prime \bullet}$ are homotopic. To see this, we define a map $h^{q}: C^{q}(\mathscr{U}, \mathscr{A}) \longrightarrow C^{q-1}(\mathscr{V}, \mathscr{A})$ by

$$
\left(h^{q} c\right)_{\beta_{0} \ldots \beta_{q-1}}=\sum_{0 \leqslant j \leqslant q-1}(-1)^{j} c_{\rho\left(\beta_{0}\right) \ldots \rho\left(\beta_{j}\right) \rho^{\prime}\left(\beta_{j}\right) \ldots \rho^{\prime}\left(\beta_{q-1}\right) \mid V_{\beta_{0} \ldots \beta_{q-1}} . . . ~ . ~}
$$

The homotopy identity $\delta^{q-1} \circ h^{q}+h^{q+1} \circ \delta^{q}=\rho^{\prime q}-\rho^{q}$ is easy to verify. Hence $\rho^{\bullet}$ and $\rho^{\prime \bullet}$ induce a map depending only on $\mathscr{U}, \mathscr{V}$ :

$$
\begin{equation*}
H^{q}\left(\rho^{\bullet}\right)=H^{q}\left(\rho^{\prime \bullet}\right): \quad \check{H}^{q}(\mathscr{U}, \mathscr{A}) \longrightarrow \check{H}^{q}(\mathscr{V}, \mathscr{A}) \tag{5.6}
\end{equation*}
$$

Now, we want to define a direct limit $\check{H}^{q}(X, \mathscr{A})$ of the groups $\check{H}^{q}(\mathscr{U}, \mathscr{A})$ by means of the refinement mappings (5.6). In order to avoid set theoretic difficulties, the coverings used in this definition will be considered as subsets of the power set $\mathscr{P}(X)$, so that the collection of all coverings becomes actually a set.
(5.7) Definition. The Čech cohomology group $\check{H}^{q}(X, \mathscr{A})$ is the direct limit

$$
\check{H}^{q}(X, \mathscr{A})=\underset{\overleftrightarrow{u}}{\lim } \quad \check{H}^{q}(\mathscr{U}, \mathscr{A})
$$

when $\because$ runs over the collection of all open coverings of $X$. Explicitly, this means that the elements of $\check{H}^{q}(X, \mathcal{A})$ are the equivalence classes in the disjoint union of the groups $\check{H}^{q}(\mathscr{U}, \mathscr{A})$, with an element in $\check{H}^{q}(\mathscr{U}, \mathscr{A})$ and another in $\check{H}^{q}(\mathscr{V}, \mathscr{A})$ identified if their images in $\check{H}^{q}(\mathscr{W}, \mathscr{A})$ coincide for some refinement $\mathscr{W}$ of the coverings $\mathscr{U}$ and $\mathscr{V}$.

If $\varphi: \mathscr{A} \rightarrow \mathscr{B}$ is a sheaf morphism, we have an obvious induced morphism $\varphi^{\bullet}$ : $C^{\bullet}(\mathscr{U}, \mathscr{A}) \longrightarrow C^{\bullet}(\mathscr{U}, \mathscr{B})$, and therefore we find a morphism

$$
H^{q}\left(\varphi^{\bullet}\right): \check{H}^{q}(\mathscr{U}, \mathscr{A}) \longrightarrow \check{H}^{q}(\mathscr{U}, \mathscr{B})
$$

Let $0 \rightarrow \mathscr{A} \rightarrow \mathscr{B} \rightarrow \mathscr{C} \rightarrow 0$ be an exact sequence of sheaves. We have an exact sequence of groups

$$
\begin{equation*}
0 \longrightarrow C^{q}(\mathscr{U}, \mathscr{A}) \longrightarrow C^{q}(\mathscr{U}, \mathscr{B}) \longrightarrow C^{q}(\mathscr{U}, \mathscr{C}) \tag{5.8}
\end{equation*}
$$

but in general the last map is not surjective, because every section in $\mathscr{C}\left(U_{\alpha_{0}, \ldots, \alpha_{q}}\right)$ need not have a lifting in $\mathscr{B}\left(U_{\alpha_{0}, \ldots, \alpha_{q}}\right)$. The image of $C^{\bullet}(\mathscr{U}, \mathscr{B})$ in $C^{\bullet}(\mathscr{U}, \mathscr{C})$ will be denoted $C_{\mathscr{B}}^{\bullet}(\mathscr{U}, \mathscr{C})$ and called the complex of liftable cochains of $\mathscr{C}$ in $\mathscr{B}$. By construction, the sequence

$$
\begin{equation*}
0 \longrightarrow C^{q}(\mathscr{U}, \mathscr{A}) \longrightarrow C^{q}(\mathscr{U}, \mathscr{B}) \longrightarrow C_{\mathscr{B}}^{q}(\mathscr{U}, \mathscr{C}) \longrightarrow 0 \tag{5.9}
\end{equation*}
$$

is exact, thus we get a corresponding long exact sequence of cohomology

$$
\begin{equation*}
\check{H}^{q}(\mathscr{U}, \mathscr{A}) \longrightarrow \check{H}^{q}(\mathscr{U}, \mathscr{B}) \longrightarrow \check{H}_{\mathscr{B}}^{q}(\mathscr{U}, \mathscr{C}) \longrightarrow \check{H}^{q+1}(\mathscr{U}, \mathscr{A}) \longrightarrow \cdots \tag{5.10}
\end{equation*}
$$

If $\mathscr{A}$ is flabby, Prop. 4.3 shows that we have $C_{\mathscr{B}}^{q}(\mathscr{U}, \mathscr{C})=C^{q}(\mathscr{U}, \mathscr{C})$, hence $\check{H}_{\mathscr{B}}^{q}(\mathscr{U}, \mathscr{C})=$ $\check{H}^{q}(\mathscr{U}, \mathscr{C})$.
(5.11) Proposition. Let $\mathscr{A}$ be a sheaf on $X$. Assume that either
a) $\mathcal{A}$ is flabby, or:
b) $X$ is paracompact and $\mathscr{A}$ is a sheaf of modules over a soft sheaf of rings $\mathscr{R}$ on $X$.

Then $\check{H}^{q}(\mathscr{U}, \mathcal{A})=0$ for every $q \geqslant 1$ and every open covering $\mathscr{U}=\left(U_{\alpha}\right)_{\alpha \in I}$ of $X$.
Proof. b) Let $\left(\psi_{\alpha}\right)_{\alpha \in I}$ be a partition of unity in $\mathscr{R}$ subordinate to $\mathscr{U}$ (Prop. 4.17). We define a map $h^{q}: C^{q}(\mathscr{U}, \mathscr{A}) \longrightarrow C^{q-1}(\cup, \mathscr{A})$ by

$$
\begin{equation*}
\left(h^{q} c\right)_{\alpha_{0} \ldots \alpha_{q-1}}=\sum_{\nu \in I} \psi_{\nu} c_{\nu \alpha_{0} \ldots \alpha_{q-1}} \tag{5.12}
\end{equation*}
$$

where $\psi_{\nu} c_{\nu \alpha_{0} \ldots \alpha_{q-1}}$ is extended by 0 on $U_{\alpha_{0} \ldots \alpha_{q-1}} \cap \complement U_{\nu}$. It is clear that

$$
\left(\delta^{q-1} h^{q} c\right)_{\alpha_{0} \ldots \alpha_{q}}=\sum_{\nu \in I} \psi_{\nu}\left(c_{\alpha_{0} \ldots \alpha_{q}}-\left(\delta^{q} c\right)_{\nu \alpha_{0} \ldots \alpha_{q}}\right),
$$

i.e. $\delta^{q-1} h^{q}+h^{q+1} \delta^{q}=$ Id. Hence $\delta^{q} c=0$ implies $\delta^{q-1} h^{q} c=c$ if $q \geqslant 1$.
a) First we show that the result is true for the sheaf $\mathscr{A}^{[0]}$. One can find a family of sets $L_{\nu} \subset U_{\nu}$ such that $\left(L_{\nu}\right)$ is a partition of $X$. If $\psi_{\nu}$ is the characteristic function of $L_{\nu}$, Formula (5.12) makes sense for any cochain $c \in C^{q}\left(\mathscr{U}, \mathscr{A}^{[0]}\right)$ because $\mathscr{A}^{[0]}$ is a module over the ring $\mathbb{Z}^{[0]}$ of germs of arbitrary functions $X \rightarrow \mathbb{Z}$. Hence $\check{H}^{q}\left(\mathscr{U}, \mathscr{A}^{[0]}\right)=0$ for $q \geqslant 1$. We shall prove this property for all flabby sheaves by induction on $q$. Consider the exact sequence

$$
0 \longrightarrow \mathscr{A} \longrightarrow \mathscr{A}^{[0]} \longrightarrow \mathscr{C} \longrightarrow 0
$$

where $\mathscr{C}=\mathscr{A}^{[0]} / \mathscr{A}$. By the remark after (5.10), we have exact sequences

$$
\begin{aligned}
& \mathscr{A}^{[0]}(X) \longrightarrow \mathscr{C}(X) \longrightarrow \check{H}^{1}(\mathscr{U}, \mathscr{A}) \longrightarrow \check{H}^{1}\left(\mathscr{U}, \mathscr{A}^{[0]}\right)=0, \\
& \check{H}^{q}(\mathscr{U}, \mathscr{C}) \longrightarrow \check{H}^{q+1}(\mathscr{U}, \mathscr{A}) \longrightarrow \check{H}^{q+1}\left(\mathscr{U}, \mathscr{A}^{[0]}\right)=0 .
\end{aligned}
$$

Then $\mathscr{A}^{[0]}(X) \longrightarrow \mathscr{C}(X)$ is surjective by Prop. 4.3, thus $\check{H}^{1}(\mathscr{U}, \mathscr{A})=0$. By 4.3 again, $\mathscr{C}$ is flabby; the induction hypothesis $\check{H}^{q}(\mathscr{U}, \mathscr{C})=0$ implies that $\check{H}^{q+1}(\mathscr{U}, \mathscr{A})=0$.

## § 5.B. Leray's Theorem for Acyclic Coverings

We first show the existence of a natural morphism from Cech cohomology to ordinary cohomology. Let $\mathscr{U}=\left(U_{\alpha}\right)_{\alpha \in I}$ be a covering of $X$. Select a map $\lambda: X \rightarrow I$ such that $x \in U_{\lambda(x)}$ for every $x \in X$. To every cochain $c \in C^{q}(\mathscr{U}, \mathscr{A})$ we associate the section $\lambda^{q} c=f \in \mathscr{A}^{[q]}(X)$ such that

$$
\begin{equation*}
f\left(x_{0}, \ldots, x_{q}\right)=c_{\lambda\left(x_{0}\right) \ldots \lambda\left(x_{q}\right)}\left(x_{q}\right) \in \mathscr{A}_{x_{q}} ; \tag{5.13}
\end{equation*}
$$

note that the right hand side is well defined as soon as

$$
x_{0} \in X, \quad x_{1} \in U_{\lambda\left(x_{0}\right)}, \quad \ldots, \quad x_{q} \in U_{\lambda\left(x_{0}\right) \ldots \lambda\left(x_{q-1}\right)} .
$$

A comparison of (2.5) and (5.13) immediately shows that the section of $\mathscr{A}^{[q+1]}(X)$ associated to $\delta^{q} c$ is

$$
\sum_{0 \leqslant j \leqslant q+1}(-1)^{j} c_{\lambda\left(x_{0}\right) \ldots \hat{\lambda\left(x_{j}\right)} \ldots \lambda\left(x_{q+1}\right)}\left(x_{q+1}\right)=\left(d^{q} f\right)\left(x_{0}, \ldots, x_{q+1}\right) .
$$

In this way we get a morphism of complexes $\lambda^{\bullet}: C^{\bullet}(\mathscr{U}, \mathscr{A}) \longrightarrow \mathscr{A}^{[\bullet]}(X)$. There is a corresponding morphism

$$
\begin{equation*}
H^{q}\left(\lambda^{\bullet}\right): \check{H}^{q}(\mathscr{U}, \mathscr{A}) \longrightarrow H^{q}(X, \mathscr{A}) \tag{5.14}
\end{equation*}
$$

If $\mathscr{V}=\left(V_{\beta}\right)_{\beta \in J}$ is a refinement of $\mathscr{U}$ such that $V_{\beta} \subset U_{\rho(\beta)}$ and $x \in V_{\mu(x)}$ for all $x, \beta$, we get a commutative diagram

\[

\]

with $\lambda=\rho \circ \mu$. In particular, (5.6) shows that the map $H^{q}\left(\lambda^{\bullet}\right)$ in (5.14) does not depend on the choice of $\lambda$ : if $\lambda^{\prime}$ is another choice, then $H^{q}\left(\lambda^{\bullet}\right)$ and $H^{q}\left(\lambda^{\bullet \bullet}\right)$ can be both factorized through the group $\check{H}^{q}(\sqrt[V]{\mathcal{V}}, \mathscr{A})$ with $V_{x}=U_{\lambda(x)} \cap U_{\lambda^{\prime}(x)}$ and $\mu=\mathrm{Id}_{X}$. By the universal property of direct limits, we get an induced morphism

$$
\begin{equation*}
\check{H}^{q}(X, \mathscr{A}) \longrightarrow H^{q}(X, \mathscr{A}) \tag{5.15}
\end{equation*}
$$

Let $0 \rightarrow \mathscr{A} \rightarrow \mathscr{B} \rightarrow \mathscr{C} \rightarrow 0$ be an exact sequence of sheaves. There is a commutative diagram

where the vertical arrows are given by the morphisms $\lambda^{\bullet}$. We obtain therefore a commutative diagram

(5.17) Theorem (Leray). Assume that

$$
H^{s}\left(U_{\alpha_{0} \ldots \alpha_{t}}, \mathscr{A}\right)=0
$$

for all indices $\alpha_{0}, \ldots, \alpha_{t}$ and $s \geqslant 1$. Then (5.14) gives an isomorphism $\check{H}^{q}(\mathscr{U}, \mathscr{A}) \simeq$ $H^{q}(X, \mathscr{A})$.

We say that the covering $\mathscr{U}$ is acyclic (with respect to $\mathscr{A}$ ) if the hypothesis of Th. 5.17 is satisfied. Leray's theorem asserts that the cohomology groups of $\mathscr{A}$ on $X$ can be
computed by means of an arbitrary acyclic covering (if such a covering exists), without using the direct limit procedure.

Proof. By induction on $q$, the result being obvious for $q=0$. Consider the exact sequence $0 \rightarrow \mathscr{A} \rightarrow \mathscr{B} \rightarrow \mathscr{C} \rightarrow 0$ with $\mathscr{B}=\mathscr{A}^{[0]}$ and $\mathscr{C}=\mathscr{A}^{[0]} / \mathscr{A}$. As $\mathscr{B}$ is acyclic, the hypothesis on $\mathscr{A}$ and the long exact sequence of cohomology imply $H^{s}\left(U_{\alpha_{0} \ldots \alpha_{t}}, \mathscr{C}\right)=0$ for $s \geqslant 1$, $t \geqslant 0$. Moreover $C_{\mathscr{B}}^{\bullet}(\mathscr{U}, \mathscr{C})=C^{\bullet}(\mathscr{U}, \mathscr{C})$ thanks to Cor. 3.6. The induction hypothesis in degree $q$ and diagram (5.16) give

$$
\begin{array}{ccccc}
\check{H}^{q}(\mathscr{U}, \mathscr{B}) & \longrightarrow & \check{H}^{q}(\mathscr{U}, \mathscr{C}) & \longrightarrow & \check{H}^{q+1}(\mathscr{U}, \mathscr{A})
\end{array}>000
$$

hence $\check{H}^{q+1}(\mathscr{U}, \mathscr{A}) \longrightarrow H^{q+1}(X, \mathscr{A})$ is also an isomorphism.
(5.18) Remark. The morphism $H^{1}\left(\lambda^{\bullet}\right): \check{H}^{1}(\mathscr{U}, \mathscr{A}) \longrightarrow H^{1}(X, \mathscr{A})$ is always injective. Indeed, we have a commutative diagram

$$
\begin{aligned}
& \check{H}^{0}(\mathscr{U}, \mathscr{B}) \longrightarrow \check{H}_{\mathscr{B}}^{0}(\mathscr{U}, \mathscr{C}) \longrightarrow \check{H}^{1}(\mathscr{U}, \mathscr{A}) \longrightarrow 0 \\
& \downarrow=~ \downarrow \downarrow \\
& H^{0}(X, \mathscr{B}) \longrightarrow H^{0}(X, \mathscr{C}) \longrightarrow H^{1}(X, \mathscr{A}) \longrightarrow 0,
\end{aligned}
$$

where $\check{H}_{\mathscr{B}}^{0}(\mathscr{U}, \mathscr{C})$ is the subspace of $\mathscr{C}(X)=H^{0}(X, \mathscr{C})$ consisting of sections which can be lifted in $\mathscr{B}$ over each $U_{\alpha}$. As a consequence, the refinement mappings

$$
H^{1}\left(\rho^{\bullet}\right): \check{H}^{1}(\mathscr{U}, \mathscr{A}) \longrightarrow \check{H}^{1}(\mathscr{V}, \mathscr{A})
$$

are also injective.

## § 5.C. Čech Cohomology on Paracompact Spaces

We will prove here that Čech cohomology theory coincides with the ordinary one on paracompact spaces.
(5.19) Proposition. Assume that $X$ is paracompact. If

$$
0 \longrightarrow \mathscr{A} \longrightarrow \mathscr{B} \longrightarrow \mathscr{C} \longrightarrow 0
$$

is an exact sequence of sheaves, there is an exact sequence

$$
\check{H}^{q}(X, \mathscr{A}) \longrightarrow \check{H}^{q}(X, \mathscr{B}) \longrightarrow \check{H}^{q}(X, \mathscr{C}) \longrightarrow \check{H}^{q+1}(X, \mathscr{A}) \longrightarrow \cdots
$$

which is the direct limit of the exact sequences (5.10) over all coverings $\mathscr{U}$.
Proof. We have to show that the natural map

$$
\xrightarrow{\lim } \check{H}_{\mathscr{B}}^{q}(\mathscr{U}, \mathscr{C}) \longrightarrow \xrightarrow{\lim } \check{H}^{q}(\mathscr{U}, \mathscr{C})
$$

is an isomorphism. This follows easily from the following lemma, which says essentially that every cochain in $\mathscr{C}$ becomes liftable in $\mathscr{B}$ after a refinement of the covering.
(5.20) Lifting lemma. Let $\mathscr{U}=\left(U_{\alpha}\right)_{\alpha \in I}$ be an open covering of $X$ and $c \in C^{q}(\mathscr{U}, \mathscr{C})$. If $X$ is paracompact, there exists a finer covering $\mathscr{V}=\left(V_{\beta}\right)_{\beta \in J}$ and a refinement map $\rho: J \rightarrow I$ such that $\rho^{q} c \in C_{\mathfrak{B}}^{q}(\mathcal{V}, \mathscr{C})$.

Proof. Since $\mathscr{U}$ admits a locally finite refinement, we may assume that $\mathscr{U}$ itself is locally finite. There exists an open covering $\mathfrak{W}=\left(W_{\alpha}\right)_{\alpha \in I}$ of $X$ such that $\bar{W}_{\alpha} \subset U_{\alpha}$. For every point $x \in X$, we can select an open neighborhood $V_{x}$ of $x$ with the following properties:
a) if $x \in W_{\alpha}$, then $V_{x} \subset W_{\alpha}$;
b) if $x \in U_{\alpha}$ or if $V_{x} \cap W_{\alpha} \neq \emptyset$, then $V_{x} \subset U_{\alpha}$;
c) if $x \in U_{\alpha_{0} \ldots \alpha_{q}}$, then $c_{\alpha_{0} \ldots \alpha_{q}} \in C^{q}\left(U_{\alpha_{0} \ldots \alpha_{q}}, \mathscr{C}^{\prime}\right)$ admits a lifting in $\mathscr{B}\left(V_{x}\right)$.

Indeed, a) (resp. c)) can be achieved because $x$ belongs to only finitely many sets $W_{\alpha}$ (resp. $U_{\alpha}$ ), and so only finitely many sections of $\mathscr{C}$ have to be lifted in $\mathscr{B}$. b) can be achieved because $x$ has a neighborhood $V_{x}^{\prime}$ that meets only finitely many sets $U_{\alpha}$; then we take

$$
V_{x} \subset V_{x}^{\prime} \cap \bigcap_{U_{\alpha} \ni x} U_{\alpha} \cap \bigcap_{U_{\alpha} \not \supset x}\left(V_{x}^{\prime} \backslash \bar{W}_{\alpha}\right) .
$$

Choose $\rho: X \rightarrow I$ such that $x \in W_{\rho(x)}$ for every $x$. Then a) implies $V_{x} \subset W_{\rho(x)}$, so $\mathscr{V}=\left(V_{x}\right)_{x \in X}$ is finer than $\mathscr{U}$, and $\rho$ defines a refinement map. If $V_{x_{0} \ldots x_{q}} \neq \emptyset$, we have

$$
V_{x_{0}} \cap W_{\rho\left(x_{j}\right)} \supset V_{x_{0}} \cap V_{x_{j}} \neq \emptyset \quad \text { for } \quad 0 \leqslant j \leqslant q,
$$

thus $V_{x_{0}} \subset U_{\rho\left(x_{0}\right) \ldots \rho\left(x_{q}\right)}$ by b). Now, c) implies that the section $c_{\rho\left(x_{0}\right) \ldots \rho\left(x_{q}\right)}$ admits a lifting in $\mathscr{B}\left(V_{x_{0}}\right)$, and in particular in $\mathscr{B}\left(V_{x_{0} \ldots x_{q}}\right)$. Therefore $\rho^{q} c$ is liftable in $\mathscr{B}$.
(5.21) Theorem. If $X$ is a paracompact topological space, the canonical morphism $H^{q}(X, \mathscr{A}) \simeq H^{q}(X, \mathscr{A})$ is an isomorphism.

Proof. Argue by induction on $q$ as in Leray's theorem, with the Čech cohomology exact sequence over $\mathscr{U}$ replaced by its direct limit in (5.16).

In the next chapters, we will be concerned only by paracompact spaces, and most often in fact by manifolds that are either compact or countable at infinity. In these cases, we will not distinguish $H^{q}(X, \mathscr{A})$ and $\check{H}^{q}(X, \mathscr{A})$.

## $\S$ 5.D. Alternate Čech Cochains

For explicit calculations, it is sometimes useful to consider a slightly modified Čech complex which has the advantage of producing much smaller cochain groups. If $\mathscr{A}$ is a sheaf and $\mathscr{U}=\left(U_{\alpha}\right)_{\alpha \in I}$ an open covering of $X$, we let $A C^{q}(\mathscr{U}, \mathscr{A}) \subset C^{q}(\mathscr{U}, \mathscr{A})$ be the subgroup of alternate Čech cochains, consisting of Čech cochains $c=\left(c_{\alpha_{0} \ldots \alpha_{q}}\right)$ such that

$$
\left\{\begin{align*}
c_{\alpha_{0} \ldots \alpha_{q}} & =0 & \text { if } \alpha_{i}=\alpha_{j}, \quad i \neq j,  \tag{5.22}\\
c_{\alpha_{\sigma(0)} \ldots \alpha_{\sigma(q)}} & =\varepsilon(\sigma) c_{\alpha_{0} \ldots \alpha_{q}} &
\end{align*}\right.
$$

for any permutation $\sigma$ of $\{1, \ldots, q\}$ of signature $\varepsilon(\sigma)$. Then the Čech differential (5.1) of an alternate cochain is still alternate, so $A C^{\bullet}(\mathscr{U}, \mathscr{A})$ is a subcomplex of $C^{\bullet}(\mathscr{U}, \mathscr{A})$. We are going to show that the inclusion induces an isomorphism in cohomology:

$$
\begin{equation*}
H^{q}\left(A C^{\bullet}(\mathscr{U}, \mathscr{A})\right) \simeq H^{q}\left(C^{\bullet}(\mathscr{U}, \mathscr{A})\right)=\check{H}^{q}(\mathscr{U}, \mathscr{A}) \tag{5.23}
\end{equation*}
$$

Select a total ordering on the index set $I$. For each such ordering, we can define a projection $\pi^{q}: C^{q}(\ddots, \mathscr{A}) \longrightarrow A C^{q}(\cup, \mathscr{A}) \subset C^{q}(\cup, \mathscr{A})$ by

$$
c \longmapsto \text { alternate } \widetilde{c} \text { such that } \widetilde{c}_{\alpha_{0} \ldots \alpha_{q}}=c_{\alpha_{0} \ldots \alpha_{q}} \text { whenever } \alpha_{0}<\ldots<\alpha_{q} .
$$

As $\pi^{\bullet}$ is a morphism of complexes, it is enough to verify that $\pi^{\bullet}$ is homotopic to the identity on $C \bullet(\mathscr{U}, \mathscr{A})$. For a given multi-index $\alpha=\left(\alpha_{0}, \ldots, \alpha_{q}\right)$, which may contain repeated indices, there is a unique permutation $(m(0), \ldots, m(q))$ of $(0, \ldots, q)$ such that

$$
\alpha_{m(0)} \leqslant \ldots \leqslant \alpha_{m(q)} \quad \text { and } m(l)<m(l+1) \text { whenever } \quad \alpha_{m(l)}=\alpha_{m(l+1)} .
$$

For $p \leqslant q$, we let $\varepsilon(\alpha, p)$ be the sign of the permutation

$$
(0, \ldots, q) \longmapsto(m(0), \ldots, m(p-1), 0,1, \ldots, \widehat{m(0)}, \ldots, m(\widehat{p-1}), \ldots, q)
$$

if the elements $\alpha_{m(0)}, \ldots, \alpha_{m(p)}$ are all distinct, and $\varepsilon(\alpha, p)=0$ otherwise. Finally, we set $h^{q}=0$ for $q \leqslant 0$ and

$$
\left(h^{q} c\right)_{\alpha_{0} \ldots \alpha_{q-1}}=\sum_{0 \leqslant p \leqslant q-1}(-1)^{p} \varepsilon(\alpha, p) c_{\alpha_{m(0)} \ldots \alpha_{m(p)} \alpha_{0} \alpha_{1} \ldots} \widehat{\alpha_{m(0)} \ldots \alpha_{m(p-1)} \ldots \alpha_{q-1}}
$$

for $q \geqslant 1$; observe that the index $\alpha_{m(p)}$ is repeated twice in the right hand side. A rather tedious calculation left to the reader shows that

$$
\left(\delta^{q-1} h^{q} c+h^{q+1} \delta^{q} c\right)_{\alpha_{0} \ldots \alpha_{q}}=c_{\alpha_{0} \ldots \alpha_{q}}-\varepsilon(\alpha, q) c_{\alpha_{m(0)} \ldots \alpha_{m(q)}}=\left(c-\pi^{q} c\right)_{\alpha_{0} \ldots \alpha_{q}} .
$$

An interesting consequence of the isomorphism (5.23) is the following:
(5.24) Proposition. Let $\mathscr{A}$ be a sheaf on a paracompact space $X$. If $X$ has arbitrarily fine open coverings or at least one acyclic open covering $\mathscr{U}=\left(U_{\alpha}\right)$ such that more than $n+1$ distinct sets $U_{\alpha_{0}}, \ldots, U_{\alpha_{n}}$ have empty intersection, then $H^{q}(X, \mathscr{A})=0$ for $q>n$.

Proof. In fact, we have $A C^{q}(\mathscr{U}, \mathscr{A})=0$ for $q>n$.

## § 6. The De Rham-Weil Isomorphism Theorem

In § 3 we defined cohomology groups by means of the simplicial flabby resolution. We show here that any resolution by acyclic sheaves could have been used instead. Let $\left(\mathscr{L}^{\bullet}, d\right)$ be a resolution of a sheaf $\mathscr{A}$. We assume in addition that all $\mathscr{L}^{q}$ are acyclic on $X$, i.e. $H^{s}\left(X, \mathscr{L}^{q}\right)=0$ for all $q \geqslant 0$ and $s \geqslant 1$. Set $\mathscr{L}^{q}=\operatorname{ker} d^{q}$. Then $\mathscr{L}^{0}=\mathscr{A}$ and for every $q \geqslant 1$ we get a short exact sequence

$$
0 \longrightarrow \mathscr{L}^{q-1} \longrightarrow \mathscr{L}^{q-1} \xrightarrow{d^{q-1}} \not \mathscr{L}^{q} \longrightarrow 0
$$

Theorem 3.5 yields an exact sequence

$$
\begin{equation*}
H^{s}\left(X, \mathscr{L}^{q-1}\right) \xrightarrow{d^{q-1}} H^{s}\left(X, \mathscr{L}^{q}\right) \xrightarrow{\partial^{s, q}} H^{s+1}\left(X, \mathscr{L}^{q-1}\right) \rightarrow H^{s+1}\left(X, \mathscr{L}^{q-1}\right)=0 . \tag{6.1}
\end{equation*}
$$

If $s \geqslant 1$, the first group is also zero and we get an isomorphism

$$
\partial^{s, q}: H^{s}\left(X, \nVdash^{q}\right) \xrightarrow{\simeq} H^{s+1}\left(X, \nVdash^{q-1}\right) .
$$

For $s=0$ we have $H^{0}\left(X, \mathscr{L}^{q-1}\right)=\mathscr{L}^{q-1}(X)$ and $H^{0}\left(X, \mathscr{L}^{q}\right)=\mathscr{L}^{q}(X)$ is the $q$-cocycle group of $\mathscr{L}^{\bullet}(X)$, so the connecting map $\partial^{0, q}$ gives an isomorphism

$$
H^{q}\left(\mathscr{L}^{\bullet}(X)\right)=\mathscr{L ^ { q }}(X) / d^{q-1} \mathscr{L}^{q-1}(x) \xrightarrow{\widetilde{\partial}^{0, q}} H^{1}\left(X, \mathscr{L}^{q-1}\right)
$$

The composite map $\partial^{q-1,1} \circ \cdots \circ \partial^{1, q-1} \circ \widetilde{\partial}^{0, q}$ therefore defines an isomorphism

$$
\begin{equation*}
H^{q}\left(\mathscr{L}^{\bullet}(X)\right) \xrightarrow{\widetilde{\partial}^{0, q}} H^{1}\left(X, \mathscr{L}^{q-1}\right) \xrightarrow{\partial^{1, q-1}} \cdots \xrightarrow{\partial^{q-1,1}} H^{q}\left(X, \mathscr{Z}^{0}\right)=H^{q}(X, \mathscr{A}) . \tag{6.2}
\end{equation*}
$$

This isomorphism behaves functorially with respect to morphisms of resolutions. Our assertion means that for every sheaf morphism $\varphi: \mathscr{A} \rightarrow \mathscr{B}$ and every morphism of resolutions $\varphi^{\bullet}: \mathscr{L}^{\bullet} \longrightarrow \mathscr{M}^{\bullet}$, there is a commutative diagram

$$
\begin{gather*}
H^{s}\left(\mathscr{L}^{\bullet}(X)\right) \longrightarrow H^{s}(X, \mathscr{A}) \\
\downarrow H^{s}\left(\varphi^{\bullet}\right) \quad \downarrow H^{s}(\varphi)  \tag{6.3}\\
H^{s}(\mathscr{M} \bullet(X)) \longrightarrow H^{s}(X, \mathscr{B}) .
\end{gather*}
$$

If $\mathscr{W}^{q}=\operatorname{ker}\left(d^{q}: \mathscr{M}^{q} \rightarrow \mathscr{M}^{q+1}\right)$, the functoriality comes from the fact that we have commutative diagrams

$$
\begin{array}{ccc}
0 \rightarrow \mathscr{L}^{q-1} \rightarrow \mathscr{L}^{q-1} \rightarrow \mathscr{L}^{q} \rightarrow 0, & H^{s}\left(X, \mathscr{\mathscr { L } ^ { q } )} \xrightarrow{\partial^{s, q}} H^{s+1}\left(X, \mathscr{\mathscr { L } ^ { q - 1 } )}\right.\right. \\
& \downarrow \varphi^{q-1} \quad \downarrow \varphi^{q-1} \downarrow \varphi^{q} & \downarrow H^{s}\left(\varphi^{q}\right) \\
0 \rightarrow W^{q-1} \rightarrow M^{q-1} \rightarrow W^{q+1}\left(\varphi^{q-1}\right) \\
0, ~ & H^{s}\left(X, W^{q}\right) \xrightarrow{\partial^{s, q}} H^{s+1}\left(X, W^{q-1}\right) .
\end{array}
$$

(6.4) De Rham-Weil isomorphism theorem. If ( $\mathscr{L}^{\bullet}, d$ ) is a resolution of $\mathscr{A}$ by sheaves $\mathscr{L}^{q}$ which are acyclic on $X$, there is a functorial isomorphism

$$
H^{q}\left(\mathscr{L}^{\bullet}(X)\right) \longrightarrow H^{q}(X, \mathscr{A})
$$

(6.5) Example: De Rham cohomology. Let $X$ be a $n$-dimensional paracompact differential manifold. Consider the resolution

$$
0 \rightarrow \mathbb{R} \rightarrow \mathscr{E}^{0} \xrightarrow{d} \mathscr{E}^{1} \rightarrow \cdots \rightarrow \mathscr{E}^{q} \xrightarrow{d} \mathscr{E}^{q+1} \rightarrow \cdots \rightarrow \mathscr{E}^{n} \rightarrow 0
$$

given by the exterior derivative $d$ acting on germs of $\mathscr{C}^{\infty}$ differential $q$-forms (cf. Example 2.2). The De Rham cohomology groups of $X$ are precisely

$$
\begin{equation*}
H_{\mathrm{DR}}^{q}(X, \mathbb{R})=H^{q}\left(\mathscr{E}^{\bullet}(X)\right) \tag{6.6}
\end{equation*}
$$

All sheaves $\mathscr{E}^{q}$ are $\mathscr{E}_{X}$-modules, so $\mathscr{E}^{q}$ is acyclic by Cor. 4.19. Therefore, we get an isomorphism

$$
\begin{equation*}
H_{\mathrm{DR}}^{q}(X, \mathbb{R}) \xrightarrow{\simeq} H^{q}(X, \mathbb{R}) \tag{6.7}
\end{equation*}
$$

from the De Rham cohomology onto the cohomology with values in the constant sheaf $\mathbb{R}$. Instead of using $\mathscr{C}^{\infty}$ differential forms, one can consider the resolution of $\mathbb{R}$ given by the exterior derivative $d$ acting on currents:

$$
0 \rightarrow \mathbb{R} \rightarrow \mathscr{D}_{n}^{\prime} \stackrel{d}{\rightarrow} \mathscr{D}_{n-1}^{\prime} \rightarrow \cdots \rightarrow \mathscr{D}_{n-q}^{\prime} \xrightarrow{d} \mathscr{D}_{n-q-1}^{\prime} \rightarrow \cdots \rightarrow \mathscr{D}_{0}^{\prime} \rightarrow 0 .
$$

The sheaves $\mathscr{D}_{q}^{\prime}$ are also $\mathscr{E}_{X}$-modules, hence acyclic. Thanks to (6.3), the inclusion $\mathscr{E}^{q} \subset \mathscr{D}_{n-q}^{\prime}$ induces an isomorphism

$$
\begin{equation*}
H^{q}\left(\mathscr{E}^{\bullet}(X)\right) \simeq H^{q}\left(\mathscr{D}_{n-\bullet}^{\prime}(X)\right), \tag{6.8}
\end{equation*}
$$

both groups being isomorphic to $H^{q}(X, \mathbb{R})$. The isomorphism between cohomology of differential forms and singular cohomology (another topological invariant) was first established by [De Rham 1931]. The above proof follows essentially the method given by [Weil 1952], in a more abstract setting. As we will see, the isomorphism (6.7) can be put under a very explicit form in terms of Čech cohomology. We need a simple lemma.
(6.9) Lemma. Let $X$ be a paracompact differentiable manifold. There are arbitrarily fine open coverings $\vartheta=\left(U_{\alpha}\right)$ such that all intersections $U_{\alpha_{0} \ldots \alpha_{q}}$ are diffeomorphic to convex sets.

Proof. Select locally finite coverings $\Omega_{j}^{\prime} \subset \subset \Omega_{j}$ of $X$ by open sets diffeomorphic to concentric euclidean balls in $\mathbb{R}^{n}$. Let us denote by $\tau_{j k}$ the transition diffeomorphism from the coordinates in $\Omega_{k}$ to those in $\Omega_{j}$. For any point $a \in \Omega_{j}^{\prime}$, the function $x \mapsto|x-a|^{2}$ computed in terms of the coordinates of $\Omega_{j}$ becomes $\left|\tau_{j k}(x)-\tau_{j k}(a)\right|^{2}$ on any patch $\Omega_{k} \ni a$. It is clear that these functions are strictly convex at $a$, thus there is a euclidean ball $B(a, \varepsilon) \subset \Omega_{j}^{\prime}$ such that all functions are strictly convex on $B(a, \varepsilon) \cap \Omega_{k}^{\prime} \subset \Omega_{k}$ (only a finite number of indices $k$ is involved). Now, choose $\mathscr{U}$ to be a (locally finite) covering of $X$ by such balls $U_{\alpha}=B\left(a_{\alpha}, \varepsilon_{\alpha}\right)$ with $U_{\alpha} \subset \Omega_{\rho(\alpha)}^{\prime}$. Then the intersection $U_{\alpha_{0} \ldots \alpha_{q}}$ is defined in $\Omega_{k}, k=\rho\left(\alpha_{0}\right)$, by the equations

$$
\left|\tau_{j k}(x)-\tau_{j k}\left(a_{\alpha_{m}}\right)\right|^{2}<\varepsilon_{\alpha_{m}}^{2}
$$

where $j=\rho\left(\alpha_{m}\right), 0 \leqslant m \leqslant q$. Hence the intersection is convex in the open coordinate chart $\Omega_{\rho\left(\alpha_{0}\right)}$.

Let $\Omega$ be an open subset of $\mathbb{R}^{n}$ which is starshaped with respect to the origin. Then the De Rham complex $\mathbb{R} \longrightarrow \mathscr{E}(\Omega)$ is acyclic: indeed, Poincaré's lemma yields a homotopy operator $k^{q}: \mathscr{E}^{q}(\Omega) \longrightarrow \mathscr{E}^{q-1}(\Omega)$ such that

$$
\begin{aligned}
& k^{q} f_{x}\left(\xi_{1}, \ldots, \xi_{q-1}\right)=\int_{0}^{1} t^{q-1} f_{t x}\left(x, \xi_{1}, \ldots, \xi_{q-1}\right) d t, \quad x \in \Omega, \quad \xi_{j} \in \mathbb{R}^{n}, \\
& k^{0} f=f(0) \in \mathbb{R} \quad \text { for } \quad f \in \mathscr{E}^{0}(\Omega)
\end{aligned}
$$

Hence $H_{\mathrm{DR}}^{q}(\Omega, \mathbb{R})=0$ for $q \geqslant 1$. Now, consider the resolution $\mathscr{E} \bullet$ of the constant sheaf $\mathbb{R}$ on $X$, and apply the proof of the De Rham-Weil isomorphism theorem to Čech cohomology groups over a covering $\mathscr{U}$ chosen as in Lemma 6.9. Since the intersections $U_{\alpha_{0} \ldots \alpha_{s}}$ are convex, all Čech cochains in $C^{s}\left(\mathscr{U}, \mathscr{\not}^{q}\right)$ are liftable in $\mathscr{E}^{q-1}$ by means of $k^{q}$. Hence
for all $s=1, \ldots, q$ we have isomorphisms $\partial^{s, q-s}: \check{H}^{s}\left(\mathscr{U}, \not \mathscr{Z}^{q-s}\right) \longrightarrow \check{H}^{s+1}\left(\mathscr{U}, \mathscr{Z}^{q-s-1}\right)$ for $s \geqslant 1$ and we get a resulting isomorphism

$$
\partial^{q-1,1} \circ \ldots \circ \partial^{1, q-1} \circ \widetilde{\partial}^{0, q}: H_{\mathrm{DR}}^{q}(X, \mathbb{R}) \xrightarrow{\simeq} \check{H}^{q}(\ddots, \mathbb{R})
$$

We are going to compute the connecting homomorphisms $\partial^{s, q-s}$ and their inverses explicitly.

Let $c$ in $C^{s}\left(\mathscr{U}, \mathscr{Z}^{q-s}\right)$ such that $\delta^{s} c=0$. As $c_{\alpha_{0} \ldots \alpha_{s}}$ is $d$-closed, we can write $c=$ $d\left(k^{q-s} c\right)$ where the cochain $k^{q-s} c \in C^{s}\left(\cup, \mathscr{E}^{q-s-1}\right)$ is defined as the family of sections $k^{q-s} c_{\alpha_{0} \ldots \alpha_{s}} \in \mathscr{E}^{q-s-1}\left(U_{\alpha_{0} \ldots \alpha_{s}}\right)$. Then $d\left(\delta^{s} k^{q-s} c\right)=\delta^{s}\left(d k^{q-s} c\right)=\delta^{s} c=0$ and

$$
\partial^{s, q-s}\{c\}=\left\{\delta^{s} k^{q-s} c\right\} \in \check{H}^{s+1}\left(\mathscr{U}, \mathscr{Z}^{q-s-1}\right) .
$$

The isomorphism $H_{\mathrm{DR}}^{q}(X, \mathbb{R}) \xrightarrow{\simeq} \check{H}^{q}(\mathscr{U}, \mathbb{R})$ is thus defined as follows: to the cohomology class $\{f\}$ of a closed $q$-form $f \in \mathscr{E}^{q}(X)$, we associate the cocycle $\left(c_{\alpha}^{0}\right)=\left(f_{\upharpoonright U_{\alpha}}\right) \in$ $C^{0}\left(\ddots, \mathscr{Z}^{q}\right)$, then the cocycle

$$
c_{\alpha \beta}^{1}=k^{q} c_{\beta}^{0}-k^{q} c_{\alpha}^{0} \in C^{1}\left(\mathscr{U}, \mathscr{\mathscr { L } ^ { q - 1 } ) , ~}\right.
$$

and by induction cocycles $\left(c_{\alpha_{0} \ldots \alpha_{s}}^{s}\right) \in C^{s}\left(\mathscr{U}, \mathscr{Z}^{q-s}\right)$ given by

$$
\begin{equation*}
c_{\alpha_{0} \ldots \alpha_{s+1}}^{s+1}=\sum_{0 \leqslant j \leqslant s+1}(-1)^{j} k^{q-s} c_{\alpha_{0} \ldots \widehat{\alpha_{j} \ldots \alpha_{s+1}}}^{s} \quad \text { on } \quad U_{\alpha_{0} \ldots \alpha_{s+1}} . \tag{6.10}
\end{equation*}
$$

The image of $\{f\}$ in $\check{H}^{q}(\mathscr{U}, \mathbb{R})$ is the class of the $q$-cocycle $\left(c_{\alpha_{0} \ldots \alpha_{q}}^{q}\right)$ in $C^{q}(\mathscr{U}, \mathbb{R})$.
Conversely, let $\left(\psi_{\alpha}\right)$ be a $\mathscr{C}^{\infty}$ partition of unity subordinate to $\mathscr{U}$. Any Čech cocycle $c \in C^{s+1}\left(\mathscr{U}, \mathscr{L}^{q-s-1}\right)$ can be written $c=\delta^{s} \gamma$ with $\gamma \in C^{s}\left(\mathscr{U}, \mathscr{E}^{q-s-1}\right)$ given by

$$
\gamma_{\alpha_{0} \ldots \alpha_{s}}=\sum_{\nu \in I} \psi_{\nu} c_{\nu \alpha_{0} \ldots \alpha_{s}},
$$

(cf. Prop. 5.11 b )), thus $\left\{c^{\prime}\right\}=\left(\partial^{s, q-s}\right)^{-1}\{c\}$ can be represented by the cochain $c^{\prime}=$ $d \gamma \in C^{s}\left(\mathscr{U}, \mathscr{O}^{q-s}\right)$ such that

$$
c_{\alpha_{0} \ldots \alpha_{s}}^{\prime}=\sum_{\nu \in I} d \psi_{\nu} \wedge c_{\nu \alpha_{0} \ldots \alpha_{s}}=(-1)^{q-s-1} \sum_{\nu \in I} c_{\nu \alpha_{0} \ldots \alpha_{s}} \wedge d \psi_{\nu} .
$$

For a reason that will become apparent later, we shall in fact modify the sign of our isomorphism $\partial^{s, q-s}$ by the factor $(-1)^{q-s-1}$. Starting from a class $\{c\} \in \check{H}^{q}(\mathscr{U}, \mathbb{R})$, we obtain inductively $\{b\} \in \check{H}^{0}\left(\mathscr{U}, \mathscr{Z}^{q}\right)$ such that

$$
\begin{equation*}
b_{\alpha_{0}}=\sum_{\nu_{0}, \ldots, \nu_{q-1}} c_{\nu_{0} \ldots \nu_{q-1} \alpha_{0}} d \psi_{\nu_{0}} \wedge \ldots \wedge d \psi_{\nu_{q-1}} \quad \text { on } \quad U_{\alpha_{0}} \tag{6.11}
\end{equation*}
$$

corresponding to $\{f\} \in H_{\mathrm{DR}}^{q}(X, \mathbb{R})$ given by the explicit formula

$$
\begin{equation*}
f=\sum_{\nu_{q}} \psi_{\nu_{q}} b_{\nu_{q}}=\sum_{\nu_{0}, \ldots, \nu_{q}} c_{\nu_{0} \ldots \nu_{q}} \psi_{\nu_{q}} d \psi_{\nu_{0}} \wedge \ldots \wedge d \psi_{\nu_{q-1}} . \tag{6.12}
\end{equation*}
$$

The choice of sign corresponds to (6.2) multiplied by $(-1)^{q(q-1) / 2}$.
(6.13) Example: Dolbeault cohomology groups. Let $X$ be a $\mathbb{C}$-analytic manifold of dimension $n$, and let $\mathscr{E}^{p, q}$ be the sheaf of germs of $\mathscr{C}^{\infty}$ differential forms of type ( $p, q$ ) with complex values. For every $p=0,1, \ldots, n$, the Dolbeault-Grothendieck Lemma I-2.9 shows that $\left(\mathscr{E}^{p, \bullet}, d^{\prime \prime}\right)$ is a resolution of the sheaf $\Omega_{X}^{p}$ of germs of holomorphic forms of degree $p$ on $X$. The Dolbeault cohomology groups of $X$ already considered in chapter 1 can be defined by

$$
\begin{equation*}
H^{p, q}(X, \mathbb{C})=H^{q}\left(\mathscr{C}^{p, \bullet}(X)\right) \tag{6.14}
\end{equation*}
$$

The sheaves $\mathscr{E}^{p, q}$ are acyclic, so we get the Dolbeault isomorphism theorem, originally proved in [Dolbeault 1953], which relates $d^{\prime \prime}$-cohomology and sheaf cohomology:

$$
\begin{equation*}
H^{p, q}(X, \mathbb{C}) \xrightarrow{\simeq} H^{q}\left(X, \Omega_{X}^{p}\right) . \tag{6.15}
\end{equation*}
$$

The case $p=0$ is especially interesting:

$$
\begin{equation*}
H^{0, q}(X, \mathbb{C}) \simeq H^{q}\left(X, \mathbb{O}_{X}\right) \tag{6.16}
\end{equation*}
$$

As in the case of De Rham cohomology, there is an inclusion $\mathscr{E}^{p, q} \subset \mathscr{D}_{n-p, n-q}^{\prime}$ and the complex of currents $\left(\mathscr{D}_{n-p, n-\bullet}^{\prime}, d^{\prime \prime}\right)$ defines also a resolution of $\Omega_{X}^{p}$. Hence there is an isomorphism:

$$
\begin{equation*}
H^{p, q}(X, \mathbb{C})=H^{q}\left(\mathscr{E}^{p, \bullet}(X)\right) \simeq H^{q}\left(\mathscr{D}_{n-p, n-\bullet}^{\prime}(X)\right) \tag{6.17}
\end{equation*}
$$

## § 7. Cohomology with Supports

As its name indicates, cohomology with supports deals with sections of sheaves having supports in prescribed closed sets. We first introduce what is an admissible family of supports.
(7.1) Definition. A family of supports on a topological space $X$ is a collection $\Phi$ of closed subsets of $X$ with the following two properties:
a) If $F, F^{\prime} \in \Phi$, then $F \cup F^{\prime} \in \Phi$;
b) If $F \in \Phi$ and $F^{\prime} \subset F$ is closed, then $F^{\prime} \in \Phi$.
(7.2) Example. Let $S$ be an arbitrary subset of $X$. Then the family of all closed subsets of $X$ contained in $S$ is a family of supports.
(7.3) Example. The collection of all compact (non necessarily Hausdorff) subsets of $X$ is a family of supports, which will be denoted simply $c$ in the sequel.
(7.4) Definition. For any sheaf $\mathscr{A}$ and any family of supports $\Phi$ on $X, \mathscr{A}_{\Phi}(X)$ will denote the set of all sections $f \in \mathscr{A}(X)$ such that $\operatorname{Supp} f \in \Phi$.

It is clear that $\mathscr{A}_{\Phi}(X)$ is a subgroup of $\mathscr{A}(X)$. We can now introduce cohomology groups with arbitrary supports.
(7.5) Definition. The cohomology groups of $\mathscr{A}$ with supports in $\Phi$ are

$$
H_{\Phi}^{q}(X, \mathscr{A})=H^{q}\left(\mathscr{A}_{\Phi}^{[\bullet]}(X)\right)
$$

The cohomology groups with compact supports will be denoted $H_{c}^{q}(X, \mathscr{A})$ and the cohomology groups with supports in a subset $S$ will be denoted $H_{S}^{q}(X, \mathbb{A})$.

In particular $H_{\Phi}^{0}(X, \mathscr{A})=\mathscr{A}_{\Phi}(X)$. If $0 \rightarrow \mathscr{A} \rightarrow \mathscr{B} \rightarrow \mathscr{C} \rightarrow 0$ is an exact sequence, there are corresponding exact sequences

$$
\begin{align*}
0 & \longrightarrow \mathscr{A}_{\Phi}^{[q]}(X)  \tag{7.6}\\
& \longrightarrow \mathscr{B}_{\Phi}^{[q]}(X) \longrightarrow \mathscr{C}_{\Phi}^{[q]}(X) \longrightarrow \cdots \\
H_{\Phi}^{q}(X, \mathscr{A}) & \longrightarrow H_{\Phi}^{q}(X, \mathscr{B}) \longrightarrow H_{\Phi}^{q}(X, \mathscr{C}) \longrightarrow H_{\Phi}^{q+1}(X, \mathscr{A}) \longrightarrow \cdots .
\end{align*}
$$

When $\mathscr{A}$ is flabby, there is an exact sequence

$$
\begin{equation*}
0 \longrightarrow \mathscr{A}_{\Phi}(X) \longrightarrow \mathscr{B}_{\Phi}(X) \longrightarrow \mathscr{C}_{\Phi}(X) \longrightarrow 0 \tag{7.7}
\end{equation*}
$$

and every $g \in \mathscr{C}_{\Phi}(X)$ can be lifted to $v \in \mathscr{B}_{\Phi}(X)$ without enlarging the support: apply the proof of Prop. 4.3 to a maximal lifting which extends $w=0$ on $W=\complement$ (Supp $g$ ). It follows that a flabby sheaf $\mathscr{A}$ is $\Phi$-acyclic, i.e. $H_{\Phi}^{q}(X, \mathscr{A})=0$ for all $q \geqslant 1$. Similarly, assume that $X$ is paracompact and that $\mathscr{A}$ is soft, and suppose that $\Phi$ has the following additional property: every set $F \in \Phi$ has a neighborhood $G \in \Phi$. An adaptation of the proofs of Prop. 4.3 and 4.13 shows that (7.7) is again exact. Therefore every soft sheaf is also $\Phi$-acyclic in that case.

As a consequence of (7.6), any resolution $\mathscr{L} \bullet$ of $\mathscr{A}$ by $\Phi$-acyclic sheaves provides a canonical De Rham-Weil isomorphism

$$
\begin{equation*}
H^{q}\left(\mathscr{L}_{\Phi}^{\bullet}(X)\right) \longrightarrow H_{\Phi}^{q}(X, \mathscr{A}) \tag{7.8}
\end{equation*}
$$

(7.9) Example: De Rham cohomology with compact support. In the special case of the De Rham resolution $\mathbb{R} \longrightarrow \mathscr{E}^{\bullet}$ on a paracompact manifold, we get an isomorphism

$$
\begin{equation*}
H_{\mathrm{DR}, c}^{q}(X, \mathbb{R}):=H^{q}\left(\left(\mathscr{D}^{\bullet}(X)\right) \xrightarrow{\simeq} H_{c}^{q}(X, \mathbb{R})\right. \tag{7.10}
\end{equation*}
$$

where $\mathscr{D}^{q}(X)$ is the space of smooth differential $q$-forms with compact support in $X$. These groups are called the De Rham cohomology groups of $X$ with compact support. When $X$ is oriented, $\operatorname{dim} X=n$, we can also consider the complex of compactly supported currents:

$$
0 \longrightarrow \mathscr{E}_{n}^{\prime}(X) \xrightarrow{d} \mathscr{E}_{n-1}^{\prime}(X) \longrightarrow \cdots \longrightarrow \mathscr{E}_{n-q}^{\prime}(X) \xrightarrow{d} \mathscr{E}_{n-q-1}^{\prime}(X) \longrightarrow \cdots
$$

Note that $\mathscr{D}^{\bullet}(X)$ and $\mathscr{E}_{n-\bullet}^{\prime}(X)$ are respectively the subgroups of compactly supported sections in $\mathscr{E}^{\bullet}$ and $\mathscr{D}_{n-\bullet}^{\prime}$, both of which are acyclic resolutions of $\mathbb{R}$. Therefore the inclusion $\mathscr{D}^{\bullet}(X) \subset \mathscr{E}_{n-\bullet}^{\prime}(X)$ induces an isomorphism

$$
\begin{equation*}
H^{q}\left(\mathscr{D}^{\bullet}(X)\right) \simeq H^{q}\left(\mathscr{E}_{n-\bullet}^{\prime}(X)\right) \tag{7.11}
\end{equation*}
$$

both groups being isomorphic to $H_{c}^{q}(X, \mathbb{R})$.

Now, we concentrate our attention on cohomology groups with compact support. We assume until the end of this section that $X$ is a locally compact space.
(7.12) Proposition. There is an isomorphism

$$
H_{c}^{q}(X, \mathscr{A})=\underset{U \subset \subset X}{\lim _{U}} H^{q}\left(\bar{U}, A_{U}\right)
$$

where $\mathscr{A}_{U}$ is the sheaf of sections of $\mathscr{A}$ vanishing on $X \backslash U$ (cf. §3).
Proof. By definition

$$
H_{c}^{q}(X, \mathscr{A})=H^{q}\left(\mathscr{A}_{c}^{[\bullet]}(X)\right)=\underset{U \subset \subset X}{\lim _{C}} H^{q}\left(\left(\mathscr{A}^{[\bullet \bullet}\right)_{U}(\bar{U})\right)
$$

since sections of $\left(\mathscr{A}^{[\bullet]}\right)_{U}(\bar{U})$ can be extended by 0 on $X \backslash U$. However, $\left(\mathscr{A}^{[\bullet]}\right)_{U}$ is a resolution of $\mathscr{A}_{U}$ and $\left(\mathscr{A}^{[q]}\right)_{U}$ is a $\mathbb{Z}^{[q]}$-module, so it is acyclic on $\bar{U}$. The De Rham-Weil isomorphism theorem implies

$$
H^{q}\left(\left(\mathscr{A}^{[\bullet]}\right)_{U}(\bar{U})\right)=H^{q}\left(\bar{U}, \mathscr{A}_{U}\right)
$$

and the proposition follows. The reader should take care of the fact that $\left(\mathscr{A}^{[q]}\right)_{U}$ does not coincide in general with $\left(\mathscr{A}_{U}\right)^{[q]}$.

The cohomology groups with compact support can also be defined by means of Čech cohomology.
(7.13) Definition. Assume that $X$ is a separable locally compact space. If $\because\left(U_{\alpha}\right)$ is a locally finite covering of $X$ by relatively compact open subsets, we let $C_{c}^{q}(\mathscr{U}, \mathcal{A})$ be the subgroups of cochains such that only finitely many coefficients $c_{\alpha_{0} \ldots \alpha_{q}}$ are non zero. The Čech cohomology groups with compact support are defined by

$$
\begin{aligned}
\check{H}_{c}^{q}(\mathscr{U}, \mathscr{A}) & =H^{q}\left(C_{c}^{\bullet}(\mathscr{U}, \mathscr{A})\right) \\
\check{H}_{c}^{q}(X, \mathscr{A}) & =\underset{\mathscr{U}}{\lim } H^{q}\left(C_{c}^{\bullet}(\mathscr{U}, \mathscr{A})\right)
\end{aligned}
$$

For such coverings $\because$, Formula (5.13) yields canonical morphisms

$$
\begin{equation*}
H^{q}\left(\lambda^{\bullet}\right): \quad \check{H}_{c}^{q}(\mathscr{U}, \mathscr{A}) \longrightarrow H_{c}^{q}(X, \mathscr{A}) . \tag{7.14}
\end{equation*}
$$

Now, the lifting Lemma 5.20 is valid for cochains with compact supports, and the same proof as the one given in $\S 5$ gives an isomorphism

$$
\begin{equation*}
\check{H}_{c}^{q}(X, \mathscr{A}) \simeq H_{c}^{q}(X, \mathscr{A}) . \tag{7.15}
\end{equation*}
$$

## § 8. Cup Product

Let $\mathscr{R}$ be a sheaf of commutative rings and $\mathscr{A}, \mathscr{B}$ sheaves of $\mathscr{R}$-modules on a space $X$. We denote by $\mathscr{A} \otimes_{\mathscr{R}} \mathscr{B}$ the sheaf on $X$ defined by

$$
\begin{equation*}
\left(\mathscr{A} \otimes_{\mathscr{R}} \mathscr{B}\right)_{x}=\mathscr{A}_{x} \otimes_{\mathscr{R}_{x}} \mathscr{B}_{x}, \tag{8.1}
\end{equation*}
$$

with the weakest topology such that the range of any section given by $\mathscr{A}(U) \otimes_{\mathscr{R}(U)} \mathscr{B}(U)$ is open in $\mathscr{A} \otimes_{\mathscr{R}} \mathscr{B}$ for any open set $U \subset X$. Given $f \in \mathscr{A}_{x}^{[p]}$ and $g \in \mathscr{B}_{x}^{[q]}$, the cup product $f \smile g \in\left(\mathscr{A} \otimes_{\mathscr{R}} \mathscr{B}\right)_{x}^{[p+q]}$ is defined by

$$
\begin{equation*}
f \smile g\left(x_{0}, \ldots, x_{p+q}\right)=f\left(x_{0}, \ldots, x_{p}\right)\left(x_{p+q}\right) \otimes g\left(x_{p}, \ldots, x_{p+q}\right) \tag{8.2}
\end{equation*}
$$

A simple computation shows that

$$
\begin{equation*}
d^{p+q}(f \smile g)=\left(d^{p} f\right) \smile g+(-1)^{p} f \smile\left(d^{q} g\right) \tag{8.3}
\end{equation*}
$$

In particular, $f \smile g$ is a cocycle if $f, g$ are cocycles, and we have

$$
\left(f+d^{p-1} f^{\prime}\right) \smile\left(g+d^{q-1} g^{\prime}\right)=f \smile g+d^{p+q-1}\left(f^{\prime} \smile g+(-1)^{p} f \smile g^{\prime}+f^{\prime} \smile d g^{\prime}\right)
$$

Consequently, there is a well defined $\mathscr{R}(X)$-bilinear morphism

$$
\begin{equation*}
H^{p}(X, \mathscr{A}) \times H^{q}(X, \mathscr{B}) \longrightarrow H^{p+q}\left(X, \mathscr{A} \otimes_{\mathscr{R}} \mathscr{B}\right) \tag{8.4}
\end{equation*}
$$

which maps a pair $(\{f\},\{g\})$ to $\{f \smile g\}$.
Let $0 \rightarrow \mathscr{B} \rightarrow \mathscr{B}^{\prime} \rightarrow \mathscr{B}^{\prime \prime} \rightarrow 0$ be an exact sequence of sheaves. Assume that the sequence obtained after taking the tensor product by $\mathscr{A}$ is also exact:

$$
0 \longrightarrow \mathscr{A} \otimes_{\mathscr{R}} \mathscr{B} \longrightarrow \mathscr{A} \otimes_{\mathscr{R}} \mathscr{B}^{\prime} \longrightarrow \mathscr{A} \otimes_{\mathscr{R}} \mathscr{B}^{\prime \prime} \longrightarrow 0
$$

Then we obtain connecting homomorphisms

$$
\begin{aligned}
\partial^{q}: & H^{q}\left(X, \mathscr{B}^{\prime \prime}\right) \longrightarrow H^{q+1}(X, \mathscr{B}), \\
\partial^{q}: & H^{q}\left(X, \mathscr{A} \otimes_{\mathscr{R}} \mathscr{B}^{\prime \prime}\right) \longrightarrow H^{q+1}\left(X, \mathscr{A} \otimes_{\mathscr{R}} \mathscr{B}\right) .
\end{aligned}
$$

For every $\alpha \in H^{p}(X, \mathscr{A}), \beta^{\prime \prime} \in H^{q}\left(X, \mathscr{B}^{\prime \prime}\right)$ we have

$$
\begin{align*}
& \partial^{p+q}\left(\alpha \smile \beta^{\prime \prime}\right)=(-1)^{p} \alpha \smile\left(\partial^{q} \beta^{\prime \prime}\right)  \tag{8.5}\\
& \partial^{p+q}\left(\beta^{\prime \prime} \smile \alpha\right)=\left(\partial^{q} \beta^{\prime \prime}\right) \smile \alpha \tag{8.5'}
\end{align*}
$$

where the second line corresponds to the tensor product of the exact sequence by $\mathscr{A}$ on the right side. These formulas are deduced from (8.3) applied to a representant $f \in \mathscr{A}^{[p]}(X)$ of $\alpha$ and to a lifting $g^{\prime} \in \mathscr{B}^{\prime[q]}(X)$ of a representative $g^{\prime \prime}$ of $\beta^{\prime \prime}$ (note that $d^{p} f=0$ ).
(8.6) Associativity and anticommutativity. Let $i: \mathscr{A} \otimes_{\mathscr{R}} \mathscr{B} \longrightarrow \mathscr{B} \otimes_{\mathscr{R}} \mathscr{A}$ be the canonical isomorphism $s \otimes t \mapsto t \otimes s$. For all $\alpha \in H^{p}(X, \mathscr{A}), \beta \in H^{q}(X, \mathscr{B})$ we have

$$
\beta \smile \alpha=(-1)^{p q} i(\alpha \smile \beta) .
$$

If $\mathscr{C}$ is another sheaf of $\mathscr{R}$-modules and $\gamma \in H^{r}(X, \mathscr{C})$, then

$$
(\alpha \smile \beta) \smile \gamma=\alpha \smile(\beta \smile \gamma)
$$

Proof. The associativity property is easily seen to hold already for all cochains

$$
(f \smile g) \smile h=f \smile(g \smile h), \quad f \in \mathscr{A}_{x}^{[p]}, \quad g \in \mathscr{B}_{x}^{[q]}, \quad h \in \mathscr{C}_{x}^{[r]} .
$$

The commutation property is obvious for $p=q=0$, and can be proved in general by induction on $p+q$. Assume for example $q \geqslant 1$. Consider the exact sequence

$$
0 \longrightarrow \mathscr{B} \longrightarrow \mathscr{B}^{\prime} \longrightarrow \mathscr{B}^{\prime \prime} \longrightarrow 0
$$

where $\mathscr{B}^{\prime}=\mathscr{B}^{[0]}$ and $\mathscr{B}^{\prime \prime}=\mathscr{B}^{[0]} / \mathscr{B}$. This exact sequence splits on each stalk (but not globally, nor even locally): a left inverse $\mathscr{B}_{x}^{[0]} \rightarrow \mathscr{B}_{x}$ of the inclusion is given by $g \mapsto g(x)$. Hence the sequence remains exact after taking the tensor product with $\mathscr{A}$. Now, as $\mathscr{B}^{\prime}$ is acyclic, the connecting homomorphism $H^{q-1}\left(X, \mathscr{B}^{\prime \prime}\right) \longrightarrow H^{q}(X, \mathscr{B})$ is onto, so there is $\beta^{\prime \prime} \in H^{q-1}\left(X, \mathscr{B}^{\prime \prime}\right)$ such that $\beta=\partial^{q-1} \beta^{\prime \prime}$. Using (8.5'), (8.5) and the induction hypothesis, we find

$$
\begin{aligned}
\beta \smile \alpha & =\partial^{p+q-1}\left(\beta^{\prime \prime} \smile \alpha\right)=\partial^{p+q-1}\left((-1)^{p(q-1)} i\left(\alpha \smile \beta^{\prime \prime}\right)\right) \\
& =(-1)^{p(q-1)} i \partial^{p+q-1}\left(\alpha \smile \beta^{\prime \prime}\right)=(-1)^{p(q-1)}(-1)^{p} i(\alpha \smile \beta) .
\end{aligned}
$$

Theorem 8.6 shows in particular that $H^{\bullet}(X, \mathscr{R})$ is a graded associative and supercommutative algebra, i.e. $\beta \smile \alpha=(-1)^{p q} \alpha \smile \beta$ for all classes $\alpha \in H^{p}(X, \mathscr{R}), \beta \in H^{q}(X, \mathscr{R})$. If $\mathscr{A}$ is a $\mathscr{R}$-module, then $H^{\bullet}(X, \mathscr{A})$ is a graded $H^{\bullet}(X, \mathscr{R})$-module.
(8.7) Remark. The cup product can also be defined for Čech cochains. Given $c \in$ $C^{p}(\mathscr{U}, \mathscr{A})$ and $c^{\prime} \in C^{q}(\mathscr{U}, \mathscr{B})$, the cochain $c \smile c^{\prime} \in C^{p+q}\left(\mathscr{U}, \mathscr{A} \otimes_{\mathscr{R}} \mathscr{B}\right)$ is defined by

$$
\left(c \smile c^{\prime}\right)_{\alpha_{0} \ldots \alpha_{p+q}}=c_{\alpha_{0} \ldots \alpha_{p}} \otimes c_{\alpha_{p} \ldots \alpha_{p+q}}^{\prime} \quad \text { on } U_{\alpha_{0} \ldots \alpha_{p+q}} .
$$

Straightforward calculations show that

$$
\delta^{p+q}\left(c \smile c^{\prime}\right)=\left(\delta^{p} c\right) \smile c^{\prime}+(-1)^{p} c \smile\left(\delta^{q} c^{\prime}\right)
$$

and that there is a commutative diagram

where the vertical arrows are the canonical morphisms $H^{s}\left(\lambda^{\bullet}\right)$ of (5.14). Note that the commutativity already holds in fact on cochains.
(8.8) Remark. Let $\Phi$ and $\Psi$ be families of supports on $X$. Then $\Phi \cap \Psi$ is again a family of supports, and Formula (8.2) defines a bilinear map

$$
\begin{equation*}
H_{\Phi}^{p}(X, \mathscr{A}) \times H_{\Psi}^{q}(X, \mathscr{B}) \longrightarrow H_{\Phi \cap \Psi}^{p+q}\left(X, \mathscr{A} \otimes_{\mathscr{R}} \mathscr{B}\right) \tag{8.9}
\end{equation*}
$$

on cohomology groups with supports. This follows immediately from the fact that $\operatorname{Supp}(f \smile g) \subset \operatorname{Supp} f \cap \operatorname{Supp} g$.
(8.10) Remark. Assume that $X$ is a differentiable manifold. Then the cohomology complex $H_{\mathrm{DR}}^{\bullet}(X, \mathbb{R})$ has a natural structure of supercommutative algebra given by the wedge product of differential forms. We shall prove the following compatibility statement:

Let $H^{q}(X, \mathbb{R}) \longrightarrow H_{\mathrm{DR}}^{q}(X, \mathbb{R})$ be the De Rham-Weil isomorphism given by Formula (6.12). Then the cup product $c^{\prime} \smile c^{\prime \prime}$ is mapped on the wedge product $f^{\prime} \wedge f^{\prime \prime}$ of the corresponding De Rham cohomology classes.

By remark 8.7, we may suppose that $c^{\prime}, c^{\prime \prime}$ are Čech cohomology classes of respective degrees $p, q$. Formulas (6.11) and (6.12) give

$$
\begin{aligned}
f_{\mid U_{\nu_{p}}}^{\prime} & =\sum_{\nu_{0}, \ldots, \nu_{p-1}} c_{\nu_{0} \ldots \nu_{p-1} \nu_{p}}^{\prime} d \psi_{\nu_{0}} \wedge \ldots \wedge d \psi_{\nu_{p-1}} \\
f^{\prime \prime} & =\sum_{\nu_{p}, \ldots, \nu_{p+q}} c_{\nu_{p} \ldots \nu_{p+q}}^{\prime \prime} \psi_{\nu_{p+q}} d \psi_{\nu_{p}} \wedge \ldots \wedge d \psi_{\nu_{p+q-1}}
\end{aligned}
$$

We get therefore

$$
f^{\prime} \wedge f^{\prime \prime}=\sum_{\nu_{0}, \ldots, \nu_{p+q}} c_{\nu_{0} \ldots \nu_{p}}^{\prime} c_{\nu_{p} \ldots \nu_{p+q}}^{\prime \prime} \psi_{\nu_{p+q}} d \psi_{\nu_{0}} \wedge \ldots \wedge \psi_{\nu_{p+q-1}},
$$

which is precisely the image of $c \smile c^{\prime}$ in the De Rham cohomology.

## § 9. Inverse Images and Cartesian Products

## § 9.A. Inverse Image of a Sheaf

Let $F: X \rightarrow Y$ be a continuous map between topological spaces $X, Y$, and let $\pi: \mathscr{A} \rightarrow Y$ be a sheaf of abelian groups. Recall that inverse image $F^{-1} \mathscr{A}$ is defined as the sheaf-space

$$
F^{-1} \mathscr{A}=\mathscr{A} \times_{Y} X=\{(s, x) ; \pi(s)=F(x)\}
$$

with projection $\pi^{\prime}=\operatorname{pr}_{2}: F^{-1} \mathscr{A} \rightarrow X$. The stalks of $F^{-1} \mathscr{A}$ are given by

$$
\begin{equation*}
\left(F^{-1} \mathscr{A}\right)_{x}=\mathscr{A}_{F(x)}, \tag{9.1}
\end{equation*}
$$

and the sections $\sigma \in F^{-1} \mathscr{A}(U)$ can be considered as continuous mappings $\sigma: U \rightarrow \mathscr{A}$ such that $\pi \circ \sigma=F$. In particular, any section $s \in \mathscr{A}(U)$ has a pull-back

$$
\begin{equation*}
F^{\star} s=s \circ F \in F^{-1} \mathscr{A}\left(F^{-1}(U)\right) . \tag{9.2}
\end{equation*}
$$

For any $v \in \mathscr{A}_{y}^{[q]}$, we define $F^{\star} v \in\left(F^{-1} \mathscr{A}\right)_{x}^{[q]}$ by

$$
\begin{equation*}
F^{\star} v\left(x_{0}, \ldots, x_{q}\right)=v\left(F\left(x_{0}\right), \ldots, F\left(x_{q}\right)\right) \in\left(F^{-1} \mathscr{A}\right)_{x_{q}}=\mathscr{A}_{F\left(x_{q}\right)} \tag{9.3}
\end{equation*}
$$

for $x_{0} \in V(x), x_{1} \in V\left(x_{0}\right), \ldots, x_{q} \in V\left(x_{0}, \ldots, x_{q-1}\right)$. We get in this way a morphism of complexes $F^{\star}: \mathscr{A}^{[\bullet]}(Y) \longrightarrow\left(F^{-1} \mathscr{A}\right)^{[\bullet]}(X)$. On cohomology groups, we thus have an induced morphism

$$
\begin{equation*}
F^{\star}: \quad H^{q}(Y, \mathscr{A}) \longrightarrow H^{q}\left(X, F^{-1} \mathscr{A}\right) \tag{9.4}
\end{equation*}
$$

Let $0 \rightarrow \mathscr{A} \rightarrow \mathscr{B} \rightarrow \mathscr{C} \rightarrow 0$ be an exact sequence of sheaves on $X$. Thanks to property (9.1), there is an exact sequence

$$
0 \longrightarrow F^{-1} \mathscr{A} \longrightarrow F^{-1} \mathscr{B} \longrightarrow F^{-1} \mathscr{C} \longrightarrow 0
$$

It is clear on the definitions that the morphism $F^{\star}$ in (9.4) commutes with the associated cohomology exact sequences. Also, $F^{\star}$ preserves the cup product, i.e. $F^{\star}(\alpha \smile \beta)=$ $F^{\star} \alpha \smile F^{\star} \beta$ whenever $\alpha, \beta$ are cohomology classes with values in sheaves $\mathscr{A}, \mathscr{B}$ on $X$. Furthermore, if $G: Y \rightarrow Z$ is a continuous map, we have

$$
\begin{equation*}
(G \circ F)^{\star}=F^{\star} \circ G^{\star} . \tag{9.5}
\end{equation*}
$$

(9.6) Remark. Similar definitions can be given for Čech cohomology. If $\mathscr{U}=\left(U_{\alpha}\right)_{\alpha \in I}$ is an open covering of $Y$, then $F^{-1} \mathscr{U}=\left(F^{-1}\left(U_{\alpha}\right)\right)_{\alpha \in I}$ is an open covering of $X$. For $c \in C^{q}(\mathscr{U}, \mathscr{A})$, we set

$$
\left(F^{\star} c\right)_{\alpha_{0} \ldots \alpha_{q}}=c_{\alpha_{0} \ldots \alpha_{q}} \circ F \in C^{q}\left(F^{-1} \mathscr{U}, F^{-1} \mathscr{A}\right)
$$

This definition is obviously compatible with the morphism from Čech cohomology to ordinary cohomology.
(9.7) Remark. Let $\Phi$ be a family of supports on $Y$. We define $F^{-1} \Psi$ to be the family of closed sets $K \subset X$ such that $F(K)$ is contained in some set $L \in \Psi$. Then (9.4) can be generalized in the form

$$
\begin{equation*}
F^{\star}: \quad H_{\Psi}^{q}(Y, \mathscr{A}) \longrightarrow H_{F^{-1} \Psi}^{q}\left(X, F^{-1} \mathscr{A}\right) . \tag{9.8}
\end{equation*}
$$

(9.9) Remark. Assume that $X$ and $Y$ are paracompact differentiable manifolds and that $F: X \rightarrow Y$ is a $\mathscr{C}^{\infty}$ map. If $\left(\psi_{\alpha}\right)_{\alpha \in I}$ is a partition of unity subordinate to $\mathscr{U}$, then $\left(\psi_{\alpha} \circ F\right)_{\alpha \in I}$ is a partition of unity on $X$ subordinate to $F^{-1} \vartheta$. Let $\left.c \in C^{q}(\vartheta), \mathbb{R}\right)$. The differential form associated to $F^{\star} c$ in the De Rham cohomology is

$$
\begin{aligned}
g & =\sum_{\nu_{0}, \ldots, \nu_{q}} c_{\nu_{0} \ldots \nu_{q}}\left(\psi_{\nu_{q}} \circ F\right) d\left(\psi_{\nu_{0}} \circ F\right) \wedge \ldots \wedge d\left(\psi_{\nu_{q-1}} \circ F\right) \\
& =F^{\star}\left(\sum_{\nu_{0}, \ldots, \nu_{q}} c_{\nu_{0} \ldots \nu_{q}} \psi_{\nu_{q}} d \psi_{\nu_{0}} \wedge \ldots \wedge d \psi_{\nu_{q-1}}\right) .
\end{aligned}
$$

Hence we have a commutative diagram

$$
\begin{aligned}
& H_{\mathrm{DR}}^{q}(X, \mathbb{R}) \xrightarrow{\simeq} \check{H}^{q}(X, \mathbb{R}) \xrightarrow{\simeq} H^{q}(X, \mathbb{R}) .
\end{aligned}
$$

## § 9.B. Cohomology Groups of a Subspace

Let $\mathscr{A}$ be a sheaf on a topological space $X$, let $S$ be a subspace of $X$ and $i_{S}: S \hookrightarrow X$ the inclusion. Then $i_{S}^{-1} \mathscr{A}$ is the restriction of $\mathscr{A}$ to $S$, so that $H^{q}(S, \mathscr{A})=H^{q}\left(S, i_{S}^{-1} \mathscr{A}\right)$ by definition. For any two subspaces $S^{\prime} \subset S$, the inclusion of $S^{\prime}$ in $S$ induces a restriction morphism

$$
H^{q}(S, \mathscr{A}) \longrightarrow H^{q}\left(S^{\prime}, \mathscr{A}\right)
$$

(9.10) Theorem. Let $\mathscr{A}$ be a sheaf on $X$ and $S$ a strongly paracompact subspace in $X$. When $\Omega$ ranges over open neighborhoods of $S$, we have

$$
H^{q}(S, \mathscr{A})=\underset{\Omega \supset S}{\lim } H^{q}(\Omega, \mathscr{A})
$$

Proof. When $q=0$, the property is equivalent to Prop. 4.7. The general case follows by induction on $q$ if we use the long cohomology exact sequences associated to the short exact sequence

$$
0 \longrightarrow \mathscr{A} \longrightarrow \mathscr{A}^{[0]} \longrightarrow \mathscr{A}^{[0]} / \mathscr{A} \longrightarrow 0
$$

on $S$ and on its neighborhoods $\Omega$ (note that the restriction of a flabby sheaf to $S$ is soft by Prop. 4.7 and the fact that every closed subspace of a strongly paracompact subspace is strongly paracompact).

## § 9.C. Cartesian Product

We use here the formalism of inverse images to deduce the cartesian product from the cup product. Let $R$ be a fixed commutative ring and $\mathscr{A} \rightarrow X, \mathscr{B} \rightarrow Y$ sheaves of $R$-modules. We define the external tensor product by

$$
\begin{equation*}
\mathscr{A} \boxtimes_{R} \mathscr{B}=\operatorname{pr}_{1}^{-1} \mathscr{A} \otimes_{R} \operatorname{pr}_{2}^{-1} \mathscr{B} \tag{9.11}
\end{equation*}
$$

where $\mathrm{pr}_{1}, \mathrm{pr}_{2}$ are the projections of $X \times Y$ onto $X, Y$ respectively. The sheaf $\mathscr{A}^{( } \boxtimes_{R} \mathscr{B}$ is thus the sheaf on $X \times Y$ whose stalks are

$$
\begin{equation*}
\left(\mathscr{A}^{\boxtimes_{R}} \mathscr{B}\right)_{(x, y)}=\mathscr{A}_{x} \otimes_{R} \mathscr{B}_{y} . \tag{9.12}
\end{equation*}
$$

For all cohomology classes $\alpha \in H^{p}(X, \mathscr{A}), \beta \in H^{q}(Y, \mathscr{B})$ the cartesian product $\alpha \times \beta \in$ $H^{p+q}\left(X \times Y, \mathscr{A} \boxtimes_{R} \mathscr{B}\right)$ is defined by

$$
\begin{equation*}
\alpha \times \beta=\left(\operatorname{pr}_{1}^{\star} \alpha\right) \smile\left(\operatorname{pr}_{2}^{\star} \beta\right) . \tag{9.13}
\end{equation*}
$$

More generally, let $\Phi$ and $\Psi$ be families of supports in $X$ and $Y$ respectively. If $\Phi \times \Psi$ denotes the family of all closed subsets of $X \times Y$ contained in products $K \times L$ of elements $K \in \Phi, L \in \Psi$, the cartesian product defines a $R$-bilinear map

$$
\begin{equation*}
H_{\Phi}^{p}(X, \mathscr{A}) \times H_{\Psi}^{q}(Y, \mathscr{B}) \longrightarrow H_{\Phi \times \Psi}^{p+q}\left(X \times Y, A^{4} \boxtimes_{R} \mathscr{B}\right) . \tag{9.14}
\end{equation*}
$$

If $\mathscr{A}^{\prime} \rightarrow X, \mathscr{B}^{\prime} \rightarrow Y$ are sheaves of abelian groups and if $\alpha^{\prime}, \beta^{\prime}$ are cohomology classes of degree $p^{\prime}, q^{\prime}$ with values in $\mathscr{A}^{\prime}, \mathscr{B}^{\prime}$, one gets easily

$$
\begin{equation*}
(\alpha \times \beta) \smile\left(\alpha^{\prime} \times \beta^{\prime}\right)=(-1)^{q p^{\prime}}\left(\alpha \smile \alpha^{\prime}\right) \times\left(\beta \smile \beta^{\prime}\right) \tag{9.15}
\end{equation*}
$$

Furthermore, if $F: X^{\prime} \rightarrow X$ and $G: Y^{\prime} \rightarrow Y$ are continuous maps, then

$$
\begin{equation*}
(F \times G)^{\star}(\alpha \times \beta)=\left(F^{\star} \alpha\right) \times\left(G^{\star} \beta\right) . \tag{9.16}
\end{equation*}
$$

## § 10. Spectral Sequence of a Filtered Complex

## §10.A. Construction of the Spectral Sequence

The theory of spectral sequences consists essentially in computing the homology groups of a differential module $(K, d)$ by "successive approximations", once a filtration $F_{p}(K)$ is given in $K$ and the cohomology groups of the graded modules $G_{p}(K)$ are known. Let us first recall some standard definitions and notations concerning filtrations.
(10.1) Definition. Let $R$ be a commutative ring. A filtration of a $R$-module $M$ is a sequence of submodules $M_{p} \subset M, p \in \mathbb{Z}$, also denoted $M_{p}=F_{p}(M)$, such that $M_{p+1} \subset$ $M_{p}$ for all $p \in \mathbb{Z}, \bigcup M_{p}=M$ and $\bigcap M_{p}=\{0\}$. The associated graded module is

$$
G(M)=\bigoplus_{p \in \mathbb{Z}} G_{p}(M), \quad G_{p}(M)=M_{p} / M_{p+1}
$$

Let $(K, d)$ be a differential module equipped with a filtration $\left(K_{p}\right)$ by differential submodules (i.e. $d K_{p} \subset K_{p}$ for every $p$ ). For any number $r \in \mathbb{N} \cup\{\infty\}$, we define $Z_{r}^{p}, B_{r}^{p} \subset G_{p}(K)=K_{p} / K_{p+1}$ by

$$
\begin{array}{ll}
Z_{r}^{p}=K_{p} \cap d^{-1} K_{p+r} \bmod K_{p+1}, & Z_{\infty}^{p}=K_{p} \cap d^{-1}\{0\} \bmod K_{p+1} \\
B_{r}^{p}=K_{p} \cap d K_{p-r+1} \quad \bmod K_{p+1}, & B_{\infty}^{p}=K_{p} \cap d K \quad \bmod K_{p+1}
\end{array}
$$

(10.3) Lemma. For every $p$ and $r$, there are inclusions

$$
\ldots \subset B_{r}^{p} \subset B_{r+1}^{p} \subset \ldots \subset B_{\infty}^{p} \subset Z_{\infty}^{p} \subset \ldots \subset Z_{r+1}^{p} \subset Z_{r}^{p} \subset \ldots
$$

and the differential d induces an isomorphism

$$
\tilde{d}: \quad Z_{r}^{p} / Z_{r+1}^{p} \longrightarrow B_{r+1}^{p+r} / B_{r}^{p+r} .
$$

Proof. It is clear that $\left(Z_{r}^{p}\right)$ decreases with $r$, that $\left(B_{r}^{p}\right)$ increases with $r$, and that $B_{\infty}^{p} \subset Z_{\infty}^{p}$. By definition

$$
\begin{aligned}
& Z_{r}^{p}=\left(K_{p} \cap d^{-1} K_{p+r}\right) /\left(K_{p+1} \cap d^{-1} K_{p+r}\right), \\
& B_{r}^{p}=\left(K_{p} \cap d K_{p-r+1}\right) /\left(K_{p+1} \cap d K_{p-r+1}\right) .
\end{aligned}
$$

The differential $d$ induces a morphism

$$
Z_{r}^{p} \longrightarrow\left(d K_{p} \cap K_{p+r}\right) /\left(d K_{p+1} \cap K_{p+r}\right)
$$

whose kernel is $\left(K_{p} \cap d^{-1}\{0\}\right) /\left(K_{p+1} \cap d^{-1}\{0\}\right)=Z_{\infty}^{p}$, whence isomorphisms

$$
\begin{aligned}
& \widehat{d}: Z_{r}^{p} / Z_{\infty}^{p} \longrightarrow\left(K_{p+r} \cap d K_{p}\right) /\left(K_{p+r} \cap d K_{p+1}\right), \\
& \widetilde{d}: Z_{r}^{p} / Z_{r+1}^{p} \longrightarrow\left(K_{p+r} \cap d K_{p}\right) /\left(K_{p+r} \cap d K_{p+1}+K_{p+r+1} \cap d K_{p}\right) .
\end{aligned}
$$

The right hand side of the last arrow can be identified to $B_{r+1}^{p+r} / B_{r}^{p+r}$, for

$$
\begin{aligned}
& B_{r}^{p+r}=\left(K_{p+r} \cap d K_{p+1}\right) /\left(K_{p+r+1} \cap d K_{p+1}\right), \\
& B_{r+1}^{p+r}=\left(K_{p+r} \cap d K_{p}\right) /\left(K_{p+r+1} \cap d K_{p}\right) .
\end{aligned}
$$

Now, for each $r \in \mathbb{N}$, we define a complex $E_{r}^{\bullet}=\bigoplus_{p \in \mathbb{Z}} E_{r}^{p}$ with a differential $d_{r}$ : $E_{r}^{p} \longrightarrow E_{r}^{p+r}$ of degree $r$ as follows: we set $E_{r}^{p}=Z_{r}^{p} / B_{r}^{p}$ and take

$$
\begin{equation*}
d_{r}: \quad Z_{r}^{p} / B_{r}^{p} \longrightarrow Z_{r}^{p} / Z_{r+1}^{p} \xrightarrow{\widetilde{d}} B_{r+1}^{p+r} / B_{r}^{p+r} \longleftrightarrow Z_{r}^{p+r} / B_{r}^{p+r} \tag{10.4}
\end{equation*}
$$

where the first arrow is the obvious projection and the third arrow the obvious inclusion. Since $d_{r}$ is induced by $d$, we actually have $d_{r} \circ d_{r}=0$; this can also be seen directly by the fact that $B_{r+1}^{p+r} \subset Z_{r+1}^{p+r}$.
(10.5) Theorem and definition. There is a canonical isomorphism $E_{r+1}^{\bullet} \simeq H^{\bullet}\left(E_{r}^{\bullet}\right)$. The sequence of differential complexes $\left(E_{r}^{\bullet}, d_{r}^{\bullet}\right)$ is called the spectral sequence of the filtered differential module ( $K, d$ ).

Proof. Since $\tilde{d}$ is an isomorphism in (10.4), we have

$$
\operatorname{ker} d_{r}=Z_{r+1}^{p} / B_{r}^{p}, \quad \operatorname{Im} d_{r}=B_{r+1}^{p+r} / B_{r}^{p+r} .
$$

Hence the image of $d_{r}: E_{r}^{p-r} \longrightarrow E_{r}^{p}$ is $B_{r+1}^{p} / B_{r}^{p}$ and

$$
H^{p}\left(E_{r}^{\bullet}\right)=\left(Z_{r+1}^{p} / B_{r}^{p}\right) /\left(B_{r+1}^{p} / B_{r}^{p}\right) \simeq Z_{r+1}^{p} / B_{r+1}^{p}=E_{r+1}^{p} .
$$

(10.6) Theorem. Consider the filtration of the homology module $H(K)$ defined by

$$
F_{p}(H(K))=\operatorname{Im}\left(H\left(K_{p}\right) \longrightarrow H(K)\right) .
$$

Then there is a canonical isomorphism

$$
E_{\infty}^{p}=G_{p}(H(K)) .
$$

Proof. Clearly $F_{p}(H(K))=\left(K_{p} \cap d^{-1}\{0\}\right) /\left(K_{p} \cap d K\right)$, whereas

$$
\begin{aligned}
& Z_{\infty}^{p}=\left(K_{p} \cap d^{-1}\{0\}\right) /\left(K_{p+1} \cap d^{-1}\{0\}\right), \quad B_{\infty}^{p}=\left(K_{p} \cap d K\right) /\left(K_{p+1} \cap d K\right), \\
& E_{\infty}^{p}=Z_{\infty}^{p} / B_{\infty}^{p}=\left(K_{p} \cap d^{-1}\{0\}\right) /\left(K_{p+1} \cap d^{-1}\{0\}+K_{p} \cap d K\right) .
\end{aligned}
$$

It follows that $E_{\infty}^{p} \simeq F_{p}(H(K)) / F_{p+1}(H(K))$.
In most applications, the differential module $K$ has a natural grading compatible with the filtration. Let us consider for example the case of a cohomology complex $K^{\bullet}=$ $\bigoplus_{l \in \mathbb{Z}} K^{l}$. The filtration $K_{p}^{\bullet}=F_{p}\left(K^{\bullet}\right)$ is said to be compatible with the differential complex structure if each $K_{p}^{\bullet}$ is a subcomplex of $K^{\bullet}$, i.e.

$$
K_{p}^{\bullet}=\bigoplus_{l \in \mathbb{Z}} K_{p}^{l}
$$

where $\left(K_{p}^{l}\right)$ is a filtration of $K^{l}$. Then we define $Z_{r}^{p, q}, B_{r}^{p, q}, E_{r}^{p, q}$ to be the sets of elements of $Z_{r}^{p}, B_{r}^{p}, E_{r}^{p}$ of total degree $p+q$. Therefore

$$
\begin{array}{lll}
Z_{r}^{p, q}=K_{p}^{p+q} \cap d^{-1} K_{p+r}^{p+q+1} \bmod K_{p+1}^{p+q}, & Z_{r}^{p}=\bigoplus Z_{r}^{p, q}, \\
B_{r}^{p, q}=K_{p}^{p+q} \cap d K_{p-r+1}^{p+q-1} \bmod K_{p+1}^{p+q}, & B_{r}^{p}=\bigoplus B_{r}^{p, q}, \\
E_{r}^{p, q}=Z_{r}^{p, q} / B_{r}^{p, q}, & E_{r}^{p}=\bigoplus E_{r}^{p, q},
\end{array}
$$

and the differential $d_{r}$ has bidegree $(r,-r+1)$, i.e.

$$
\begin{equation*}
d_{r}: \quad E_{r}^{p, q} \longrightarrow E_{r}^{p+r, q-r+1} . \tag{10.8}
\end{equation*}
$$

For an element of pure bidegree $(p, q), p$ is called the filtering degree, $q$ the complementary degree and $p+q$ the total degree.
(10.9) Definition. A filtration $\left(K_{p}^{\bullet}\right)$ of a complex $K^{\bullet}$ is said to be regular if for each degree $l$ there are indices $\nu(l) \leqslant N(l)$ such that $K_{p}^{l}=K^{l}$ for $p<\nu(l)$ and $K_{p}^{l}=0$ for $p>N(l)$.

If the filtration is regular, then (10.7) and (10.7 ) show that

$$
\begin{aligned}
Z_{r}^{p, q} & =Z_{r+1}^{p, q}=\ldots=Z_{\infty}^{p, q} \quad \text { for } \quad r>N(p+q+1)-p, \\
B_{r}^{p, q} & =B_{r+1}^{p, q}=\ldots=B_{\infty}^{p, q} \quad \text { for } \quad r>p+1-\nu(p+q-1),
\end{aligned}
$$

therefore $E_{r}^{p, q}=E_{\infty}^{p, q}$ for $r \geqslant r_{0}(p, q)$. We say that the spectral sequence converges to its limit term, and we write symbolically

$$
\begin{equation*}
E_{r}^{p, q} \Longrightarrow H^{p+q}\left(K^{\bullet}\right) \tag{10.10}
\end{equation*}
$$

to express the following facts: there is a spectral sequence whose terms of the $r$-th generation are $E_{r}^{p, q}$, the sequence converges to a limit term $E_{\infty}^{p, q}$, and $E_{\infty}^{p, l-p}$ is the term $G_{p}\left(H^{l}\left(K^{\bullet}\right)\right)$ in the graded module associated to some filtration of $H^{l}\left(K^{\bullet}\right)$.
(10.11) Definition. The spectral sequence is said to collapse in $E_{r}^{\bullet}$ if all terms $Z_{k}^{p, q}$, $B_{k}^{p, q}, E_{k}^{p, q}$ are constant for $k \geqslant r$, or equivalently if $d_{k}=0$ in all bidegrees for $k \geqslant r$.
(10.12) Special case. Assume that there exists an integer $r \geqslant 2$ and an index $q_{0}$ such that $E_{r}^{p, q}=0$ for $q \neq q_{0}$. Then this property remains true for larger values of $r$, and we must have $d_{r}=0$. It follows that the spectral sequence collapses in $E_{r}^{\bullet}$ and that

$$
H^{l}\left(K^{\bullet}\right)=E_{r}^{l-q_{0}, q_{0}} .
$$

Similarly, if $E_{r}^{p, q}=0$ for $p \neq p_{0}$ and some $r \geqslant 1$ then

$$
H^{l}\left(K^{\bullet}\right)=E_{r}^{p_{0}, l-p_{0}} .
$$

## § 10.B. Computation of the First Terms

Consider an arbitrary spectral sequence. For $r=0$, we have $Z_{0}^{p}=K_{p} / K_{p+1}, B_{0}^{p}=$ $\{0\}$, thus

$$
\begin{equation*}
E_{0}^{p}=K_{p} / K_{p+1}=G_{p}(K) \tag{10.13}
\end{equation*}
$$

The differential $d_{0}$ is induced by $d$ on the quotients, and

$$
\begin{equation*}
E_{1}^{p}=H\left(G_{p}(K)\right) . \tag{10.14}
\end{equation*}
$$

Now, there is a short exact sequence of differential modules

$$
0 \longrightarrow G_{p+1}(K) \longrightarrow K_{p} / K_{p+2} \longrightarrow G_{p}(K) \longrightarrow 0
$$

We get therefore a connecting homomorphism

$$
\begin{equation*}
E_{1}^{p}=H\left(G_{p}(K)\right) \xrightarrow{\partial} H\left(G_{p+1}(K)\right)=E_{1}^{p+1} . \tag{10.15}
\end{equation*}
$$

We claim that $\partial$ coincides with the differential $d_{1}$ : indeed, both are induced by $d$. When $K^{\bullet}$ is a filtered cohomology complex, $d_{1}$ is the connecting homomorphism

$$
\begin{equation*}
E_{1}^{p, q}=H^{p+q}\left(G_{p}\left(K^{\bullet}\right)\right) \xrightarrow{\partial} H^{p+q+1}\left(G_{p+1}\left(K^{\bullet}\right)\right)=E_{1}^{p+1, q} . \tag{10.16}
\end{equation*}
$$

## §11. Spectral Sequence of a Double Complex

A double complex is a bigraded module $K^{\bullet \bullet}=\bigoplus K^{p, q}$ together with a differential $d=d^{\prime}+d^{\prime \prime}$ such that

$$
\begin{equation*}
d^{\prime}: K^{p, q} \longrightarrow K^{p+1, q}, \quad d^{\prime \prime}: K^{p, q+1} \longrightarrow K^{p, q+1} \tag{11.1}
\end{equation*}
$$

and $d \circ d=0$. In particular, $d^{\prime}$ and $d^{\prime \prime}$ satisfy the relations

$$
\begin{equation*}
d^{\prime 2}=d^{\prime \prime 2}=0, \quad d^{\prime} d^{\prime \prime}+d^{\prime \prime} d^{\prime}=0 \tag{11.2}
\end{equation*}
$$

The simple complex associated to $K^{\bullet \bullet}$ is defined by

$$
K^{l}=\bigoplus_{p+q=l} K^{p, q}
$$

together with the differential $d$. We will suppose here that both graduations of $K^{\bullet \bullet \bullet}$ are positive, i.e. $K^{p, q}=0$ for $p<0$ or $q<0$. The first filtration of $K^{\bullet}$ is defined by

$$
\begin{equation*}
K_{p}^{l}=\bigoplus_{i+j=l, i \geqslant p} K^{i, j}=\bigoplus_{p \leqslant i \leqslant l} K^{i, l-i} . \tag{11.3}
\end{equation*}
$$

The graded module associated to this filtration is of course $G_{p}\left(K^{l}\right)=K^{p, l-p}$, and the differential induced by $d$ on the quotient coincides with $d^{\prime \prime}$ because $d^{\prime}$ takes $K_{p}^{l}$ to $K_{p+1}^{l+1}$. Thus we have a spectral sequence beginning by

$$
\begin{equation*}
E_{0}^{p, q}=K^{p, q}, \quad d_{0}=d^{\prime \prime}, \quad E_{1}^{p, q}=H_{d^{\prime \prime}}^{q}\left(K^{p, \bullet}\right) \tag{11.4}
\end{equation*}
$$

By (10.16), $d_{1}$ is the connecting homomorphism associated to the short exact sequence

$$
0 \longrightarrow K^{p+1, \bullet} \longrightarrow K^{p, \bullet} \oplus K^{p+1, \bullet} \longrightarrow K^{p, \bullet} \longrightarrow 0
$$

where the differential is given by $\left(d \bmod K^{p+2, \bullet}\right)$ for the central term and by $d^{\prime \prime}$ for the two others. The definition of the connecting homomorphism in the proof of Th. 1.5 shows that

$$
d_{1}=\partial: \quad H_{d^{\prime \prime}}^{q}\left(K^{p, \bullet}\right) \longrightarrow H_{d^{\prime \prime}}^{q}\left(K^{p+1, \bullet}\right)
$$

is induced by $d^{\prime}$. Consequently, we find

$$
\begin{equation*}
E_{2}^{p, q}=H_{d^{\prime}}^{p}\left(E_{1}^{\bullet, q}\right)=H_{d^{\prime}}^{p}\left(H_{d^{\prime \prime}}^{q}\left(K^{\bullet, \bullet}\right)\right) \tag{11.5}
\end{equation*}
$$

For such a spectral sequence, several interesting additional features can be pointed out. For all $r$ and $l$, there is an injective homomorphism

$$
E_{r+1}^{0, l} \longleftrightarrow E_{r}^{0, l}
$$

whose image can be identified with the set of $d_{r}$-cocycles in $E_{r}^{0, l}$; the coboundary group is zero because $E_{r}^{p, q}=0$ for $q<0$. Similarly, $E_{r}^{l, 0}$ is equal to its cocycle submodule, and there is a surjective homomorphism

$$
E_{r}^{l, 0} \longrightarrow E_{r+1}^{l, 0} \simeq E_{r}^{l, 0} / d_{r} E_{r}^{l-r, r-1} .
$$

Furthermore, the filtration on $H^{l}\left(K^{\bullet}\right)$ begins at $p=0$ and stops at $p=l$, i.e.

$$
\begin{equation*}
F_{0}\left(H^{l}\left(K^{\bullet}\right)\right)=H^{l}\left(K^{\bullet}\right), \quad F_{p}\left(H^{l}\left(K^{\bullet}\right)\right)=0 \quad \text { for } p>l \tag{11.6}
\end{equation*}
$$

Therefore, there are canonical maps

$$
\begin{align*}
& H^{l}\left(K^{\bullet}\right) \longrightarrow G_{0}\left(H^{l}\left(K^{\bullet}\right)\right)=E_{\infty}^{0, l} \hookrightarrow E_{r}^{0, l} \\
& E_{r}^{l, 0} \longrightarrow E_{\infty}^{l, 0}=G_{l}\left(H^{l}\left(K^{\bullet}\right)\right) \longleftrightarrow H^{l}\left(K^{\bullet}\right) . \tag{11.7}
\end{align*}
$$

These maps are called the edge homomorphisms of the spectral sequence.
(11.8) Theorem. There is an exact sequence

$$
0 \longrightarrow E_{2}^{1,0} \longrightarrow H^{1}\left(K^{\bullet}\right) \longrightarrow E_{2}^{0,1} \xrightarrow{d_{2}} E_{2}^{2,0} \longrightarrow H^{2}\left(K^{\bullet}\right)
$$

where the non indicated arrows are edge homomorphisms.
Proof. By 11.6, the graded module associated to $H^{1}\left(K^{\bullet}\right)$ has only two components, and we have an exact sequence

$$
0 \longrightarrow E_{\infty}^{1,0} \longrightarrow H^{1}\left(K^{\bullet}\right) \longrightarrow E_{\infty}^{0,1} \longrightarrow 0
$$

However $E_{\infty}^{1,0}=E_{2}^{1,0}$ because all differentials $d_{r}$ starting from $E_{r}^{1,0}$ or abuting to $E_{r}^{1,0}$ must be zero for $r \geqslant 2$. Similarly, $E_{\infty}^{0,1}=E_{3}^{0,1}$ and $E_{\infty}^{2,0}=E_{3}^{2,0}$, thus there is an exact sequence

$$
0 \longrightarrow E_{\infty}^{0,1} \longrightarrow E_{2}^{0,1} \xrightarrow{d_{2}} E_{2}^{2,0} \longrightarrow E_{\infty}^{2,0} \longrightarrow 0
$$

A combination of the two above exact sequences yields

$$
0 \longrightarrow E_{2}^{1,0} \longrightarrow H^{1}\left(K^{\bullet}\right) \longrightarrow E_{2}^{0,1} \xrightarrow{d_{2}} E_{2}^{2,0} \longrightarrow E_{\infty}^{2,0} \longrightarrow 0 .
$$

Taking into account the injection $E_{\infty}^{2,0} \hookrightarrow H^{2}\left(K^{\bullet}\right)$ in (11.7), we get the required exact sequence.
(11.9) Example. Let $X$ be a complex manifold of dimension $n$. Consider the double complex $K^{p, q}=\mathscr{C}_{p, q}^{\infty}(X, \mathbb{C})$ together with the exterior derivative $d=d^{\prime}+d^{\prime \prime}$. Then there is a spectral sequence which starts from the Dolbeault cohomology groups

$$
E_{1}^{p, q}=H^{p, q}(X, \mathbb{C})
$$

and which converges to the graded module associated to a filtration of the De Rham cohomology groups:

$$
E_{r}^{p, q} \Longrightarrow H_{\mathrm{DR}}^{p+q}(X, \mathbb{C})
$$

This spectral sequence is called the Hodge-Frölicher spectral sequence [Frölicher 1955]. We will study it in much more detail in chapter 6 when $X$ is compact.

Frequently, the spectral sequence under consideration can be obtained from two distinct double complexes and one needs to compare the final cohomology groups. The following lemma can often be applied.
(11.10) Lemma. Let $K^{p, q} \longrightarrow L^{p, q}$ be a morphism of double complexes (i.e. a double sequence of maps commuting with $d^{\prime}$ and $\left.d^{\prime \prime}\right)$. Then there are induced morphisms

$$
{ }_{K} E_{r}^{\bullet, \bullet} \longrightarrow{ }_{L} E_{r}^{\bullet, \bullet}, \quad r \geqslant 0
$$

of the associated spectral sequences. If one of these morphisms is an isomorphism for some $r$, then $H^{l}\left(K^{\bullet}\right) \longrightarrow H^{l}\left(L^{\bullet}\right)$ is an isomorphism.

Proof. If the $r$-terms are isomorphic, they have the same cohomology groups, thus ${ }_{K} E_{r+1}^{\bullet \bullet \bullet} \simeq{ }_{L} E_{r+1}^{\bullet, \bullet}$ and ${ }_{K} E_{\infty}^{\bullet \bullet \bullet} \simeq{ }_{L} E_{\infty}^{\bullet \bullet \bullet}$ in the limit. The lemma follows from the fact that if a morphism of graded modules $\varphi: M \longrightarrow M^{\prime}$ induces an isomorphism $G_{\bullet}(M) \longrightarrow$ $G_{\bullet}\left(M^{\prime}\right)$, then $\varphi$ is an isomorphism.

## § 12. Hypercohomology Groups

Let $\left(\mathscr{L}^{\bullet}, \delta\right)$ be a complex of sheaves

$$
0 \longrightarrow \mathscr{L}^{0} \quad \stackrel{\delta^{0}}{\longrightarrow} \mathscr{L}^{1} \longrightarrow \cdots \longrightarrow \mathscr{L}^{q} \xrightarrow{\delta^{q}} \cdots
$$

on a topological space $X$. We denote by $\mathscr{H}^{q}=\mathscr{H}^{q}\left(\mathscr{L}^{\bullet}\right)$ the $q$-th sheaf of cohomology of this complex; thus $\mathscr{H}^{q}$ is a sheaf of abelian groups over $X$. Our goal is to define "generalized cohomology groups" attached to $\mathscr{L} \bullet$ on $X$, in such a way that these groups only depend on the cohomology sheaves $\mathscr{H}^{q}$. For this, we associate to $\mathscr{L}$ • the double complex of groups

$$
\begin{equation*}
K_{\mathscr{L}}^{p, q}=\left(\mathscr{L}^{q}\right)^{[p]}(X) \tag{12.1}
\end{equation*}
$$

with differential $d^{\prime}=d^{p}$ given by (2.5), and with $d^{\prime \prime}=(-1)^{p}\left(\delta^{q}\right)^{[p]}$. As $\left(\delta^{q}\right)^{[\bullet]}$ : $\left(\mathscr{L}^{q}\right)^{[\bullet]} \longrightarrow\left(\mathscr{L}^{q+1}\right)^{[\bullet]}$ is a morphism of complexes, we get the expected relation $d^{\prime} d^{\prime \prime}+$ $d^{\prime \prime} d^{\prime}=0$.
(12.2) Definition. The groups $H^{q}\left(K_{\dot{\mathscr{L}}}^{\bullet}\right)$ are called the hypercohomology groups of $\mathscr{L} \bullet$ and are denoted $\mathbb{H}^{q}\left(X, \mathscr{L}^{\bullet}\right)$.

Clearly $\mathbb{H}^{0}\left(X, \mathscr{L}^{\bullet}\right)=\mathscr{H}^{0}(X)$ where $\mathscr{H}^{0}=\operatorname{ker} \delta^{0}$ is the first cohomology sheaf of $\mathscr{L}^{\bullet}$. If $\varphi^{\bullet}: \mathscr{L}^{\bullet} \longrightarrow \mathscr{M}^{\bullet}$ is a morphism of sheaf complexes, there is an associated morphism of double complexes $\varphi^{\bullet \bullet \bullet}: K_{\mathscr{L}}^{\bullet \bullet \bullet} \longrightarrow K_{\mathscr{M}}^{\bullet \bullet \bullet}$, hence a natural morphism

$$
\begin{equation*}
\mathbb{H}^{q}\left(\varphi^{\bullet}\right): \mathbb{H}^{q}\left(X, \mathscr{L}^{\bullet}\right) \longrightarrow \mathbb{H}^{q}\left(X, \mathcal{M}^{\bullet}\right) \tag{12.3}
\end{equation*}
$$

We first list a few immediate properties of hypercohomology groups, whose proofs are left to the reader.
(12.4) Proposition. The following properties hold:
a) If $\mathscr{L}^{q}=0$ for $q \neq 0$, then $\mathbb{H}^{q}\left(X, \mathscr{L}^{\bullet}\right)=H^{q}\left(X, \mathscr{L}^{0}\right)$.
b) If $\mathscr{L} \bullet[s]$ denotes the complex $\mathscr{L} \bullet$ shifted of $s$ indices to the right, i.e. $\mathscr{L} \bullet[s]^{q}=\mathscr{L}^{q-s}$, then $\mathbb{H}^{q}\left(X, \mathscr{L}^{\bullet}[s]\right)=\mathbb{H}^{q-s}\left(X, \mathscr{L}^{\bullet}\right)$.
c) If $0 \longrightarrow \mathscr{L}^{\bullet} \longrightarrow \mathcal{M}^{\bullet} \longrightarrow \mathcal{N}^{\bullet} \longrightarrow 0$ is an exact sequence of sheaf complexes, there is a long exact sequence

$$
\cdots \mathbb{H}^{q}\left(X, \mathscr{L}^{\bullet}\right) \longrightarrow \mathbb{H}^{q}\left(X, \mathscr{M}^{\bullet}\right) \longrightarrow \mathbb{H}^{q}\left(X, \mathcal{N}^{\bullet}\right) \xrightarrow{\partial} \mathbb{H}^{q+1}\left(X, \mathscr{L}^{\bullet}\right) \cdots .
$$

If $\mathscr{L} \bullet$ is a sheaf complex, the spectral sequence associated to the first filtration of $K_{\mathscr{L}}^{\bullet}$ is given by

$$
E_{1}^{p, q}=H_{d^{\prime \prime}}^{q}\left(K_{\mathscr{L}}^{p, \bullet}\right)=H^{q}\left(\left(\mathscr{L}^{\bullet}\right)^{[p]}(X)\right) .
$$

However by (2.9) the functor $\mathcal{A} \longmapsto \mathscr{A}^{[p]}(X)$ preserves exact sequences. Therefore, we get

$$
\begin{align*}
& E_{1}^{p, q}=\left(\mathscr{H}^{q}\left(\mathscr{L}^{\bullet}\right)\right)^{[p]}(X),  \tag{12.5}\\
& E_{2}^{p, q}=H^{p}\left(X, \mathscr{H}^{q}\left(\mathscr{L}^{\bullet}\right)\right),
\end{align*}
$$

since $E_{2}^{p, q}=H_{d^{\prime}}^{p}\left(E_{1}^{\bullet, q}\right)$. If $\varphi^{\bullet}: \mathscr{L}^{\bullet} \longrightarrow \mathcal{M}^{\bullet}$ is a morphism, an application of Lemma 11.10 to the $E_{2}$-term of the associated first spectral sequences of $K_{\dot{\mathscr{L}}}^{\bullet \bullet \bullet}$ and $K_{\mathcal{M}}^{\bullet \bullet \bullet}$ yields:
(12.6) Corollary. If $\varphi^{\bullet}: \mathscr{L}^{\bullet} \longrightarrow \mathcal{M}^{\bullet}$ is a quasi-isomorphism (this means that $\varphi^{\bullet}$ induces an isomorphism $\left.\mathscr{H}^{\bullet}\left(\mathscr{L}^{\bullet}\right) \longrightarrow \mathscr{H}^{\bullet}\left(\mathscr{M}^{\bullet}\right)\right)$, then

$$
\mathbb{H}^{l}\left(\varphi^{\bullet}\right): \mathbb{H}^{l}\left(X, \mathscr{L}^{\bullet}\right) \longrightarrow \mathbb{H}^{l}\left(X, \mathscr{M}^{\bullet}\right)
$$

is an isomorphism.

Now, we may reverse the roles of the indices $p, q$ and of the differentials $d^{\prime}, d^{\prime \prime}$. The second filtration $F_{p}\left(K_{\mathscr{L}}^{l}\right)=\bigoplus_{j \geqslant p} K_{\mathscr{L}}^{l-j, j}$ is associated to a spectral sequence such that $\widetilde{E}_{1}^{p, q}=H_{d^{\prime}}^{q}\left(K_{\mathscr{L}}^{\bullet, p}\right)=H_{d^{\prime}}^{q}\left(\left(\mathscr{L}^{p}\right)^{[\bullet]}(X)\right)$, hence

$$
\begin{align*}
& \widetilde{E}_{1}^{p, q}=H^{q}\left(X, \mathscr{L}^{p}\right),  \tag{12.7}\\
& \widetilde{E}_{2}^{p, q}=H_{\delta}^{p}\left(H^{q}\left(X, \mathscr{L}^{\bullet}\right)\right) .
\end{align*}
$$

These two spectral sequences converge to limit terms which are the graded modules associated to filtrations of $\mathbb{H}^{\bullet}\left(X, \mathscr{L}^{\bullet}\right)$; these filtrations are in general different. Let us mention two interesting special cases.

- Assume first that the complex $\mathscr{L}^{\bullet}$ is a resolution of a sheaf $\mathscr{A}$, so that $\mathscr{H}^{0}=\mathscr{A}$ and $\mathscr{H}^{q}=0$ for $q \geqslant 1$. Then we find

$$
E_{2}^{p, 0}=H^{p}(X, \mathscr{A}), \quad E_{2}^{p, q}=0 \text { for } q \geqslant 1 .
$$

It follows that the first spectral sequence collapses in $E_{2}^{\bullet}$, and 10.12 implies

$$
\begin{equation*}
\mathbb{H}^{l}\left(X, \mathscr{L}^{\bullet}\right) \simeq H^{l}(X, \mathscr{A}) . \tag{12.8}
\end{equation*}
$$

- Now, assume that the sheaves $\mathscr{L}^{q}$ are acyclic. The second spectral sequence gives

$$
\begin{gather*}
\widetilde{E}_{2}^{p, 0}=H^{p}\left(\mathscr{L}^{\bullet}(X)\right), \quad \widetilde{E}_{2}^{p, q}=0 \text { for } q \geqslant 1, \\
\mathbb{H}^{l}\left(X, \mathscr{L}^{\bullet}\right) \simeq H^{l}\left(\mathscr{L}^{\bullet}(X)\right) . \tag{12.9}
\end{gather*}
$$

If both conditions hold, i.e. if $\mathscr{L}$ • is a resolution of a sheaf $\mathscr{A}$ by acyclic sheaves, then (12.8) and (12.9) can be combined to obtain a new proof of the De Rham-Weil isomorphism $H^{l}(X, \mathscr{A}) \simeq H^{l}\left(\mathscr{L}^{\bullet}(X)\right)$.

## § 13. Direct Images and the Leray Spectral Sequence

## $\S$ 13.A. Direct Images of a Sheaf

Let $X, Y$ be topological spaces, $F: X \rightarrow Y$ a continuous map and $\mathscr{A}$ a sheaf of abelian groups on $X$. Recall that the direct image $F_{\star} \mathscr{A}$ is the presheaf on $Y$ defined for any open set $U \subset Y$ by

$$
\begin{equation*}
\left(F_{\star} \mathscr{A}\right)(U)=\mathscr{A}\left(F^{-1}(U)\right) . \tag{13.1}
\end{equation*}
$$

Axioms (II-2.4' and (II-2.4") are clearly satisfied, thus $F_{\star} \mathscr{A}$ is in fact a sheaf. The following result is obvious:

$$
\begin{equation*}
\mathscr{A} \text { is flabby } \quad \Longrightarrow \quad F_{\star} \mathscr{A} \text { is flabby. } \tag{13.2}
\end{equation*}
$$

Every sheaf morphism $\varphi: \mathscr{A} \rightarrow \mathscr{B}$ induces a corresponding morphism

$$
F_{\star} \varphi: \quad F_{\star} \mathscr{A} \longrightarrow F_{\star} \mathscr{B},
$$

so $F_{\star}$ is a functor on the category of sheaves on $X$ to the category of sheaves on $Y$. This functor is exact on the left: indeed, to every exact sequence $0 \rightarrow \mathscr{A} \rightarrow \mathscr{B} \rightarrow \mathscr{C}$ is associated an exact sequence

$$
0 \longrightarrow F_{\star} \mathscr{A} \longrightarrow F_{\star} \mathscr{B} \longrightarrow F_{\star} \mathscr{C},
$$

but $F_{\star} \mathscr{B} \rightarrow F_{\star} \mathscr{C}$ need not be onto if $\mathscr{B} \rightarrow \mathscr{C}$ is. All this follows immediately from the considerations of $\S 3$. In particular, the simplicial flabby resolution ( $\mathscr{A}^{[\bullet]}, d$ ) yields a complex of sheaves

$$
\begin{equation*}
0 \longrightarrow F_{\star} \mathscr{A}^{[0]} \longrightarrow F_{\star} \mathscr{A}^{[1]} \longrightarrow \cdots \longrightarrow F_{\star} \mathscr{A}^{[q]} \xrightarrow{F_{\star} d^{q}} F_{\star} \mathscr{A}^{[q+1]} \longrightarrow \cdots . \tag{13.3}
\end{equation*}
$$

(13.4) Definition. The $q$-th direct image of $\mathscr{A}$ by $F$ is the $q$-th cohomology sheaf of the sheaf complex (13.3). It is denoted

$$
R^{q} F_{\star} \mathscr{A}=\mathscr{H}^{q}\left(F_{\star} \mathscr{A}^{[\bullet]}\right) .
$$

As $F_{\star}$ is exact on the left, the sequence $0 \rightarrow F_{\star} \mathscr{A} \rightarrow F_{\star} \not \mathscr{A}^{[0]} \rightarrow F_{\star} \mathscr{A}^{[1]}$ is exact, thus

$$
\begin{equation*}
R^{0} F_{\star} \mathscr{A}=F_{\star} \mathscr{A} . \tag{13.5}
\end{equation*}
$$

We now compute the stalks of $R^{q} F_{\star} \mathscr{A}$. As the kernel or cokernel of a sheaf morphism is obtained stalk by stalk, we have

$$
\left(R^{q} F_{\star} \mathscr{A}\right)_{y}=H^{q}\left(\left(F_{\star} \mathscr{A}^{[\bullet]}\right)_{y}\right)=\underset{U \ni y}{\lim _{U}} H^{q}\left(F_{\star} \mathscr{A}^{[\bullet]}(U)\right) .
$$

The very definition of $F_{\star}$ and of sheaf cohomology groups implies

$$
H^{q}\left(F_{\star} \mathscr{A}^{[\bullet]}(U)\right)=H^{q}\left(\mathscr{A}^{[\bullet]}\left(F^{-1}(U)\right)\right)=H^{q}\left(F^{-1}(U), \mathscr{A}\right),
$$

hence we find

$$
\begin{equation*}
\left(R^{q} F_{\star} \mathscr{A}\right)_{y}=\underset{U \ni y}{\lim } H^{q}\left(F^{-1}(U), \mathscr{A}\right), \tag{13.6}
\end{equation*}
$$

i.e. $R^{q} F_{\star} \mathscr{A}$ is the sheaf associated to the presheaf $U \mapsto H^{q}\left(F^{-1}(U), \mathscr{A}\right)$. We must stress here that the stronger relation $R^{q} F_{\star} \mathscr{A}(U)=H^{q}\left(F^{-1}(U), \mathscr{A}\right)$ need not be true in general. If the fiber $F^{-1}(y)$ is strongly paracompact in $X$ and if the family of open sets $F^{-1}(U)$ is a fundamental family of neighborhoods of $F^{-1}(y)$ (this situation occurs for example if $X$ and $Y$ are locally compact spaces and $F$ a proper map, or if $F=\operatorname{pr}_{1}: X=Y \times S \longrightarrow Y$ where $S$ is compact), Th. 9.10 implies the more natural relation

$$
\left(R^{q} F_{\star} \mathscr{A}\right)_{y}=H^{q}\left(F^{-1}(y), \mathscr{A}\right) .
$$

Let $0 \rightarrow \mathscr{A} \rightarrow \mathscr{B} \rightarrow \mathscr{C} \rightarrow 0$ be an exact sequence of sheaves on $X$. Apply the long exact sequence of cohomology on every open set $F^{-1}(U)$ and take the direct limit over $U$. We get an exact sequence of sheaves:

$$
\begin{align*}
0 & \longrightarrow F_{\star} \mathscr{A} \tag{13.7}
\end{align*} \longrightarrow F_{\star} \mathscr{B} \longrightarrow F_{\star} \mathscr{C} \longrightarrow R^{1} F_{\star} \mathscr{A} \longrightarrow \cdots .
$$

## § 13.B. Leray Spectral Sequence

For any continuous map $F: X \rightarrow Y$, the Leray spectral sequence relates the cohomology groups of a sheaf $\mathscr{A}$ on $X$ and those of its direct images $R^{q} F_{\star} \mathscr{A}$ on $Y$. Consider the two spectral sequences $E_{r}^{\bullet}, \widetilde{E}_{r}^{\bullet}$ associated with the complex of sheaves $\mathscr{L}^{\bullet}=F_{\star} \mathscr{A}^{[\bullet]}$ on $Y$, as in $\S 12$. By definition we have $\mathscr{H}^{q}\left(\mathscr{L}^{\bullet}\right)=R^{q} F_{\star} \mathscr{A}$. By (12.5') the second term of the first spectral sequence is

$$
E_{2}^{p, q}=H^{p}\left(Y, R^{q} F_{\star} \mathscr{A}\right),
$$

and this spectral sequence converges to the graded module associated to a filtration of $\mathbb{H}^{l}\left(Y, F_{\star} \not \mathcal{A}^{[\bullet]}\right)$. On the other hand, (13.2) implies that $F_{\star} \mathscr{A}^{[q]}$ is flabby. Hence, the second special case (12.9) can be applied:

$$
\mathbb{H}^{l}\left(Y, F_{\star} \mathscr{A}^{[\bullet]}\right) \simeq H^{l}\left(F_{\star} \mathscr{A}^{[\bullet]}(Y)\right)=H^{l}\left(\mathscr{A}^{[\bullet]}(X)\right)=H^{l}(X, \mathscr{A}) .
$$

We may conclude this discussion by the following
(13.8) Theorem. For any continuous map $F: X \rightarrow Y$ and any sheaf $\mathscr{A}$ of abelian groups on $X$, there exists a spectral sequence whose $E_{2}^{\bullet}$ term is

$$
E_{2}^{p, q}=H^{p}\left(Y, R^{q} F_{\star} \mathscr{A}\right),
$$

which converges to a limit term $E_{\infty}^{p, l-p}$ equal to the graded module associated with a filtration of $H^{l}(X, \mathscr{A})$. The edge homomorphism

$$
H^{l}\left(Y, F_{\star} \mathscr{A}\right) \longrightarrow E_{\infty}^{l, 0} \hookrightarrow H^{l}(X, \mathscr{A})
$$

coincides with the composite morphism

$$
F^{\#}: \quad H^{l}\left(Y, F_{\star} \mathscr{A}\right) \xrightarrow{F^{\star}} H^{l}\left(X, F^{-1} F_{\star} \mathscr{A}\right) \xrightarrow{H^{l}\left(\mu_{F}\right)} H^{l}(X, \mathscr{A})
$$

where $\mu_{F}: F^{-1} F_{\star} \mathscr{A} \longrightarrow \mathscr{A}$ is the canonical sheaf morphism.
Proof. Only the last statement remains to be proved. The morphism $\mu_{F}$ is defined as follows: every element $s \in\left(F^{-1} F_{\star} \mathscr{A}\right)_{x}=\left(F_{\star} \mathscr{A}\right)_{F(x)}$ is the germ of a section $\widetilde{s} \in$ $F_{\star} \mathscr{A}(V)=\mathscr{A}\left(F^{-1}(V)\right)$ on a neighborhood $V$ of $F(x)$. Then $F^{-1}(V)$ is a neighborhood of $x$ and we let $\mu_{F} s$ be the germ of $\widetilde{s}$ at $x$.

Now, we observe that to any commutative diagram of topological spaces and continuous maps is associated a commutative diagram involving the direct image sheaves and their cohomology groups:

$$
\begin{array}{cc}
X \xrightarrow{F} \underset{G}{Y} \underset{H}{ } \downarrow_{H} & H^{l}(X, \mathscr{A}) \stackrel{F^{\#}}{\rightleftarrows} H^{l}\left(Y, F_{\star} \mathscr{A}\right) \\
X^{\prime} \xrightarrow{F^{\prime}} Y^{\prime} & G^{\#} \uparrow \\
H^{l}\left(X^{\prime}, G_{\star} \mathscr{A}\right) \stackrel{F^{\prime \#}}{\rightleftarrows} H^{l}\left(Y^{\prime}, F_{\star}^{\prime} G_{\star} \mathscr{A}\right) .
\end{array}
$$

There is a similar commutative diagram in which $F^{\#}$ and $F^{\prime \#}$ are replaced by the edge homomorphisms of the spectral sequences of $F$ and $F^{\prime}$ : indeed there is a natural morphism $H^{-1} F_{\star}^{\prime} \mathscr{B} \longrightarrow F_{\star} G^{-1} \mathscr{B}$ for any sheaf $\mathscr{B}$ on $X^{\prime}$, so we get a morphism of sheaf complexes

$$
H^{-1} F_{\star}^{\prime}\left(G_{\star} \mathscr{A}\right)^{[\bullet]} \longrightarrow F_{\star} G^{-1}\left(G_{\star} \mathscr{A}\right)^{[\bullet]} \longrightarrow F_{\star}\left(G^{-1} G_{\star} \mathscr{A}\right)^{[\bullet]} \longrightarrow F_{\star} \mathscr{A}^{[\bullet]}
$$

hence also a morphism of the spectral sequences associated to both ends.
The special case $X^{\prime}=Y^{\prime}=Y, G=F, F^{\prime}=H=\operatorname{Id}_{Y}$ then shows that our statement is true for $F$ if it is true for $F^{\prime}$. Hence we may assume that $F$ is the identity map; in this case, we need only show that the edge homomorphism of the spectral sequence of
$F_{\star} \mathscr{A}^{[\bullet]}=\mathscr{A}^{[\bullet]}$ is the identity map. This is an immediate consequence of the fact that we have a quasi-isomorphism

$$
(\cdots \rightarrow 0 \rightarrow \mathscr{A} \rightarrow 0 \rightarrow \cdots) \longrightarrow \mathscr{A}^{[\bullet]} .
$$

(13.9) Corollary. If $R^{q} F_{\star} \mathscr{A}=0$ for $q \geqslant 1$, there is an isomorphism $H^{l}\left(Y, F_{\star} \mathscr{A}\right) \simeq$ $H^{l}(X, \mathscr{A})$ induced by $F^{\#}$.

Proof. We are in the special case 10.12 with $E_{2}^{p, q}=0$ for $q \neq 0$, so

$$
H^{l}\left(Y, F_{\star} \mathscr{A}\right)=E_{2}^{l, 0} \simeq H^{l}(X, \mathscr{A})
$$

(13.10) Corollary. Let $F: X \longrightarrow Y$ be a proper finite-to-one map. For any sheaf $\mathscr{A}$ on $X$, we have $R^{q} F_{\star} \mathscr{A}=0$ for $q \geqslant 1$ and there is an isomorphism $H^{l}\left(Y, F_{\star} \mathscr{A}\right) \simeq H^{l}(X, \mathscr{A})$.

Proof. By definition of higher direct images, we have

$$
\left(R^{q} F_{\star} \mathscr{A}\right)_{y}=\underset{U \ni y}{\lim } H^{q}\left(\mathscr{A}^{[\bullet]}\left(F^{-1}(U)\right)\right) .
$$

If $F^{-1}(y)=\left\{x_{1}, \ldots, x_{m}\right\}$, the assumptions imply that $\left(F^{-1}(U)\right)$ is a fundamental system of neighborhoods of $\left\{x_{1}, \ldots, x_{m}\right\}$. Therefore

$$
\left(R^{q} F_{\star} \mathscr{A}\right)_{y}=\bigoplus_{1 \leqslant j \leqslant m} H^{q}\left(\mathscr{A}_{x_{j}}^{[\bullet \bullet}\right)= \begin{cases}\bigoplus \mathscr{A}_{x_{j}} & \text { for } q=0, \\ 0 & \text { for } q \geqslant 1,\end{cases}
$$

and we conclude by Cor. 13.9.
Corollary 13.10 can be applied in particular to the inclusion $j: S \rightarrow X$ of a closed subspace $S$. Then $j_{\star} \mathscr{A}$ coincides with the sheaf $\mathscr{A}^{S}$ defined in $\S 3$ and we get $H^{q}(S, \mathscr{A})=$ $H^{q}\left(X, \mathscr{A}^{S}\right)$. It is very important to observe that the property $R^{q} j_{\star} \mathscr{A}=0$ for $q \geqslant 1$ need not be true if $S$ is not closed.

## $\S$ 13.C. Topological Dimension

As a first application of the Leray spectral sequence, we are going to derive some properties related to the concept of topological dimension.
(13.11) Definition. A non empty space $X$ is said to be of topological dimension $\leqslant n$ if $H^{q}(X, \mathscr{A})=0$ for any $q>n$ and any sheaf $\mathscr{A}$ on $X$. We let $\operatorname{topdim} X$ be the smallest such integer $n$ if it exists, and $+\infty$ otherwise.
(13.12) Criterion. For a paracompact space $X$, the following conditions are equivalent:
a) $\operatorname{top} \operatorname{dim} X \leqslant n$;
b) the sheaf $\not \mathscr{L}^{n}=\operatorname{ker}\left(\mathscr{A}^{[n]} \longrightarrow \mathscr{A}^{[n+1]}\right)$ is soft for every sheaf $\mathscr{A}$;
c) every sheaf $\mathscr{A}$ admits a resolution $0 \rightarrow \mathscr{L}^{0} \rightarrow \cdots \rightarrow \mathscr{L}^{n} \rightarrow 0$ of length $n$ by soft sheaves.

Proof. b) $\Longrightarrow$ c). Take $\mathscr{L}^{q}=\mathscr{A}^{[q]}$ for $q<n$ and $\mathscr{L}^{n}=\mathscr{L}^{n}$.
c) $\Longrightarrow$ a). For every sheaf $\mathscr{A}$, the De Rham-Weil isomorphism implies $H^{q}(X, \mathscr{A})=$ $H^{q}(\mathscr{L} \bullet(X))=0$ when $q>n$.
a) $\Longrightarrow \mathrm{b})$. Let $S$ be a closed set and $U=X \backslash S$. As in Prop. 7.12, $\left(\mathscr{A}^{[\bullet]}\right)_{U}$ is an acyclic resolution of $\mathscr{A}_{U}$. Clearly $\operatorname{ker}\left(\left(\mathscr{A}^{[n]}\right)_{U} \rightarrow\left(\mathscr{A}^{[n+1]}\right)_{U}\right)=\mathscr{L}_{U}^{n}$, so the isomorphisms (6.2) obtained in the proof of the De Rham-Weil theorem imply

$$
H^{1}\left(X, \mathscr{Z}_{U}^{n}\right) \simeq H^{n+1}\left(X, \mathscr{A}_{U}\right)=0 .
$$

By (3.10), the restriction map $\mathscr{\not}^{n}(X) \longrightarrow \mathscr{\not}^{n}(S)$ is onto, so $\mathscr{\not}^{n}$ is soft.
(13.13) Theorem. The following properties hold:
a) If $X$ is paracompact and if every point of $X$ has a neighborhood of topological dimension $\leqslant n$, then topdim $X \leqslant n$.
b) If $S \subset X$, then topdim $S \leqslant \operatorname{topdim} X$ provided that $S$ is closed or $X$ metrizable.
c) If $X, Y$ are metrizable spaces, one of them locally compact, then

$$
\operatorname{topdim}(X \times Y) \leqslant \operatorname{topdim} X+\operatorname{topdim} Y
$$

d) If $X$ is metrizable and locally homeomorphic to a subspace of $\mathbb{R}^{n}$, then topdim $X \leqslant n$.

Proof. a) Apply criterion 13.12 b) and the fact that softness is a local property (Prop. 4.12).
b) When $S$ is closed in $X$, the property follows from Cor. 13.10 . When $X$ is metrizable, any subset $S$ is strongly paracompact. Let $j: S \longrightarrow X$ be the injection and $\mathscr{A}$ a sheaf on $S$. As $\mathscr{A}=\left(j_{\star} \mathscr{A}\right)_{\mid S}$, we have

$$
H^{q}(S, \mathscr{A})=H^{q}\left(S, j_{\star} \mathscr{A}\right)=\underset{\Omega \supset S}{\lim _{\longrightarrow}} H^{q}\left(\Omega, j_{\star} \mathscr{A}\right)
$$

by Th. 9.10 . We may therefore assume that $S$ is open in $X$. Then every point of $S$ has a neighborhood which is closed in $X$, so we conclude by a) and the first case of b).
c) Thanks to a) and b), we may assume for example that $X$ is compact. Let $\mathscr{A}$ be a sheaf on $X \times Y$ and $\pi: X \times Y \longrightarrow Y$ the second projection. Set $n_{X}=\operatorname{topdim} X$, $n_{Y}=\operatorname{topdim} Y$. In virtue of (13.6'), we have $R^{q} \pi_{\star} \mathscr{A}=0$ for $q>n_{X}$. In the Leray spectral sequence, we obtain therefore

$$
E_{2}^{p, q}=H^{p}\left(Y, R^{q} \pi_{\star} \mathscr{A}\right)=0 \quad \text { for } p>n_{Y} \quad \text { or } \quad q>n_{X},
$$

thus $E_{\infty}^{p, l-p}=0$ when $l>n_{X}+n_{Y}$ and we infer $H^{l}(X \times Y, \mathscr{A})=0$.
d) The unit interval $[0,1] \subset \mathbb{R}$ is of topological dimension $\leqslant 1$, because $[0,1]$ admits arbitrarily fine coverings

$$
\begin{equation*}
\mathscr{U}_{k}=([0,1] \cap](\alpha-1) / k,(\alpha+1) / k[)_{0 \leqslant \alpha \leqslant k} \tag{13.14}
\end{equation*}
$$

for which only consecutive open sets $U_{\alpha}, U_{\alpha+1}$ intersect; we may therefore apply Prop. 5.24. Hence $\left.\mathbb{R}^{n} \simeq\right] 0,1\left[{ }^{n} \subset[0,1]^{n}\right.$ is of topological dimension $\leqslant n$ by b) and c). Property d) follows.

## § 14. Alexander-Spanier Cohomology

## $\S$ 14.A. Invariance by Homotopy

Alexander-Spanier's theory can be viewed as the special case of sheaf cohomology theory with constant coefficients, i.e. with values in constant sheaves.
(14.1) Definition. Let $X$ be a topological space, $R$ a commutative ring and $M$ a $R$ module. The constant sheaf $X \times M$ is denoted $M$ for simplicity. The Alexander-Spanier $q$-th cohomology group with values in $M$ is the sheaf cohomology group $H^{q}(X, M)$.

In particular $H^{0}(X, M)$ is the set of locally constant functions $X \rightarrow M$, therefore $H^{0}(X, M) \simeq M^{E}$, where $E$ is the set of connected components of $X$. We will not repeat here the properties of Alexander-Spanier cohomology groups that are formal consequences of those of general sheaf theory, but we focus our attention instead on new features, such as invariance by homotopy.
(14.2) Lemma. Let $I$ denote the interval $[0,1]$ of real numbers. Then $H^{0}(I, M)=M$ and $H^{q}(I, M)=0$ for $q \neq 0$.

Proof. Let us employ alternate Čech cochains for the coverings $\mathscr{U}_{n}$ defined in (13.14). As $I$ is paracompact, we have

$$
H^{q}(I, M)=\underset{\longrightarrow}{\lim } \check{H}^{q}\left(\ddots_{n}, M\right) .
$$

However, the alternate Čech complex has only two non zero components and one non zero differential:

$$
\begin{aligned}
& A C^{0}\left(U_{n}, M\right)=\left\{\left(c_{\alpha}\right)_{0 \leqslant \alpha \leqslant n}\right\}=M^{n+1}, \\
& A C^{1}\left(U_{n}, M\right)=\left\{\left(c_{\alpha \alpha+1}\right)_{0 \leqslant \alpha \leqslant n-1}\right\}=M^{n}, \\
& d^{0}:\left(c_{\alpha}\right) \longmapsto\left(c_{\alpha \alpha+1}^{\prime}\right)=\left(c_{\alpha+1}-c_{\alpha}\right) .
\end{aligned}
$$

We see that $d^{0}$ is surjective, and that $\operatorname{ker} d^{0}=\{(m, m, \ldots, m)\}=M$.
For any continuous map $f: X \longrightarrow Y$, the inverse image of the constant sheaf $M$ on $Y$ is $f^{-1} M=M$. We get therefore a morphism

$$
\begin{equation*}
f^{\star}: H^{q}(Y, M) \longrightarrow H^{q}(X, M) \tag{14.3}
\end{equation*}
$$

which, as already mentioned in $\S 9$, is compatible with cup product.
(14.4) Proposition. For any space $X$, the projection $\pi: X \times I \longrightarrow X$ and the injections $i_{t}: X \longrightarrow X \times I, x \longmapsto(x, t)$ induce inverse isomorphisms

$$
H^{q}(X, M) \underset{i_{t}^{\star}}{\stackrel{\pi^{\star}}{\rightleftarrows}} H^{q}(X \times I, M)
$$

In particular, $i_{t}^{\star}$ does not depend on $t$.
Proof. As $\pi \circ i_{t}=\mathrm{Id}$, we have $i_{t}^{\star} \circ \pi^{\star}=\mathrm{Id}$, so it is sufficient to check that $\pi^{\star}$ is an isomorphism. However $\left(R^{q} \pi_{\star} M\right)_{x}=H^{q}(I, M)$ in virtue of $\left(13.6^{\prime}\right)$, so we get

$$
R^{0} \pi_{\star} M=M, \quad R^{q} \pi_{\star} M=0 \quad \text { for } \quad q \neq 0
$$

and conclude by Cor. 13.9.
(14.4) Theorem. If $f, g: X \longrightarrow Y$ are homotopic maps, then

$$
f^{\star}=g^{\star}: H^{q}(Y, M) \longrightarrow H^{q}(X, M) .
$$

Proof. Let $H: X \times I \longrightarrow Y$ be a homotopy between $f$ and $g$, with $f=H \circ i_{0}$ and $g=H \circ i_{1}$. Proposition 14.3 implies

$$
f^{\star}=i_{0}^{\star} \circ H^{\star}=i_{1}^{\star} \circ H^{\star}=g^{\star} .
$$

We denote $f \sim g$ the homotopy equivalence relation. Two spaces $X, Y$ are said to be homotopically equivalent $(X \sim Y)$ if there exist continuous maps $u: X \longrightarrow Y$, $v: Y \longrightarrow X$ such that $v \circ u \sim \operatorname{Id}_{X}$ and $u \circ v \sim \operatorname{Id}_{Y}$. Then $H^{q}(X, M) \simeq H^{q}(Y, M)$ and $u^{\star}, v^{\star}$ are inverse isomorphisms.
(14.5) Example. A subspace $S \subset X$ is said to be a (strong) deformation retract of $X$ if there exists a retraction of $X$ onto $S$, i.e. a map $r: X \longrightarrow S$ such that $r \circ j=\operatorname{Id}_{S}$ ( $j=$ inclusion of $S$ in $X$ ), which is a deformation of $\operatorname{Id}_{X}$, i.e. there exists a homotopy $H: X \times I \longrightarrow X$ relative to $S$ between $\operatorname{Id}_{X}$ and $j \circ r:$

$$
H(x, 0)=x, \quad H(x, 1)=r(x) \quad \text { on } \quad X, \quad H(x, t)=x \quad \text { on } S \times I .
$$

Then $X$ and $S$ are homotopically equivalent. In particular $X$ is said to be contractible if $X$ has a deformation retraction onto a point $x_{0}$. In this case

$$
H^{q}(X, M)=H^{q}\left(\left\{x_{0}\right\}, M\right)= \begin{cases}M & \text { for } q=0 \\ 0 & \text { for } q \neq 0 .\end{cases}
$$

(14.6) Corollary. If $X$ is a compact differentiable manifold, the cohomology groups $H^{q}(X, R)$ are finitely generated over $R$.

Proof. Lemma 6.9 shows that $X$ has a finite covering $\mathscr{U}$ such that the intersections $U_{\alpha_{0} \ldots \alpha_{q}}$ are contractible. Hence $\mathscr{U}$ is acyclic, $H^{q}(X, R)=H^{q}\left(C^{\bullet}(U, R)\right)$ and each Čech cochain space is a finitely generated (free) module.
(14.7) Example: Cohomology Groups of Spheres. Set

$$
S^{n}=\left\{x \in \mathbb{R}^{n+1} ; x_{0}^{2}+x_{1}^{2}+\ldots+x_{n}^{2}=1\right\}, \quad n \geqslant 1 .
$$

We will prove by induction on $n$ that

$$
H^{q}\left(S^{n}, M\right)= \begin{cases}M & \text { for } q=0 \text { or } q=n  \tag{14.8}\\ 0 & \text { otherwise }\end{cases}
$$

As $S^{n}$ is connected, we have $H^{0}\left(S^{n}, M\right)=M$. In order to compute the higher cohomology groups, we apply the Mayer-Vietoris exact sequence 3.11 to the covering ( $U_{1}, U_{2}$ ) with

$$
U_{1}=S^{n} \backslash\{(-1,0, \ldots, 0)\}, \quad U_{2}=S^{n} \backslash\{(1,0, \ldots, 0)\} .
$$

Then $U_{1}, U_{2} \approx \mathbb{R}^{n}$ are contractible, and $U_{1} \cap U_{2}$ can be retracted by deformation on the equator $S^{n} \cap\left\{x_{0}=0\right\} \approx S^{n-1}$. Omitting $M$ in the notations of cohomology groups, we get exact sequences

$$
\begin{align*}
& H^{0}\left(U_{1}\right) \oplus H^{0}\left(U_{2}\right) \longrightarrow H^{0}\left(U_{1} \cap U_{2}\right) \longrightarrow H^{1}\left(S^{n}\right) \longrightarrow 0, \\
& 0 \longrightarrow H^{q-1}\left(U_{1} \cap U_{2}\right) \longrightarrow H^{q}\left(S^{n}\right) \longrightarrow 0, \quad q \geqslant 2 \tag{14.9"}
\end{align*}
$$

For $n=1, U_{1} \cap U_{2}$ consists of two open arcs, so we have

$$
H^{0}\left(U_{1}\right) \oplus H^{0}\left(U_{2}\right)=H^{0}\left(U_{1} \cap U_{2}\right)=M \times M
$$

and the first arrow in $\left(14.9^{\prime}\right)$ is $\left(m_{1}, m_{2}\right) \longmapsto\left(m_{2}-m_{1}, m_{2}-m_{1}\right)$. We infer easily that $H^{1}\left(S^{1}\right)=M$ and that

$$
H^{q}\left(S^{1}\right)=H^{q-1}\left(U_{1} \cap U_{2}\right)=0 \quad \text { for } \quad q \geqslant 2
$$

For $n \geqslant 2, U_{1} \cap U_{2}$ is connected, so the first arrow in (14.9') is onto and $H^{1}\left(S^{n}\right)=0$. The second sequence (14.9") yields $H^{q}\left(S^{n}\right) \simeq H^{q-1}\left(S^{n-1}\right)$. An induction concludes the proof.

## §14.B. Relative Cohomology Groups and Excision Theorem

Let $X$ be a topological space and $S$ a subspace. We denote by $M^{[q]}(X, S)$ the subgroup of sections $u \in M^{[q]}(X)$ such that $u\left(x_{0}, \ldots, x_{q}\right)=0$ when

$$
\left(x_{0}, \ldots, x_{q}\right) \in S^{q}, \quad x_{1} \in V\left(x_{0}\right), \ldots, x_{q} \in V\left(x_{0}, \ldots, x_{q-1}\right)
$$

Then $M^{[\bullet]}(X, S)$ is a subcomplex of $M^{[\bullet]}(X)$ and we define the relative cohomology group of the pair $(X, S)$ with values in $M$ as

$$
\begin{equation*}
H^{q}(X, S ; M)=H^{q}\left(M^{[\bullet]}(X, S)\right) \tag{14.10}
\end{equation*}
$$

By definition of $M^{[q]}(X, S)$, there is an exact sequence

$$
\begin{equation*}
0 \longrightarrow M^{[q]}(X, S) \longrightarrow M^{[q]}(X) \longrightarrow\left(M_{\upharpoonright S}\right)^{[q]}(S) \longrightarrow 0 \tag{14.11}
\end{equation*}
$$

The reader should take care of the fact that $\left(M_{\upharpoonright S}\right)^{[q]}(S)$ does not coincide with the module of sections $M^{[q]}(S)$ of the sheaf $M^{[q]}$ on $X$, except if $S$ is open. The snake lemma shows that there is an "exact sequence of the pair":

$$
\begin{equation*}
H^{q}(X, S ; M) \rightarrow H^{q}(X, M) \rightarrow H^{q}(S, M) \rightarrow H^{q+1}(X, S ; M) \cdots \tag{14.12}
\end{equation*}
$$

We have in particular $H^{0}(X, S ; M)=M^{E}$, where $E$ is the set of connected components of $X$ which do not meet $S$. More generally, for a triple $(X, S, T)$ with $X \supset S \supset T$, there is an "exact sequence of the triple":

$$
\begin{gather*}
0 \longrightarrow M^{[q]}(X, S) \longrightarrow M^{[q]}(X, T) \longrightarrow M^{[q]}(S, T) \longrightarrow 0 \\
H^{q}(X, S ; M) \longrightarrow H^{q}(X, T ; M) \longrightarrow H^{q}(S, T ; M) \longrightarrow H^{q+1}(X, S ; M)
\end{gather*}
$$

The definition of the cup product in (8.2) shows that $\alpha \smile \beta$ vanishes on $S \cup S^{\prime}$ if $\alpha, \beta$ vanish on $S, S^{\prime}$ respectively. Therefore, we obtain a bilinear map

$$
\begin{equation*}
H^{p}(X, S ; M) \times H^{q}\left(X, S^{\prime} ; M^{\prime}\right) \longrightarrow H^{p+q}\left(X, S \cup S^{\prime} ; M \otimes M^{\prime}\right) \tag{14.13}
\end{equation*}
$$

If $f:(X, S) \longrightarrow(Y, T)$ is a morphism of pairs, i.e. a continuous map $X \rightarrow Y$ such that $f(S) \subset T$, there is an induced pull-back morphism

$$
\begin{equation*}
f^{\star}: \quad H^{q}(Y, T ; M) \longrightarrow H^{q}(X, S ; M) \tag{14.14}
\end{equation*}
$$

which is compatible with the cup product. Two morphisms of pairs $f, g$ are said to be homotopic when there is a pair homotopy $H:(X \times I, S \times I) \longrightarrow(Y, T)$. An application of the exact sequence of the pair shows that

$$
\pi^{\star}: H^{q}(X, S ; M) \longrightarrow H^{q}(X \times I, S \times I ; M)
$$

is an isomorphism. It follows as above that $f^{\star}=g^{\star}$ as soon as $f, g$ are homotopic.
(14.15) Excision theorem. For subspaces $\bar{T} \subset S^{\circ}$ of $X$, the restriction morphism $H^{q}(X, S ; M) \longrightarrow H^{q}(X \backslash T, S \backslash T ; M)$ is an isomorphism.

Proof. Under our assumption, it is not difficult to check that the surjective restriction $\operatorname{map} M^{[q]}(X, S) \longrightarrow M^{[q]}(X \backslash T, S \backslash T)$ is also injective, because the kernel consists of sections $u \in M^{[q]}(X)$ such that $u\left(x_{0}, \ldots, x_{q}\right)=0$ on $(X \backslash T)^{q+1} \cup S^{q+1}$, and this set is a neighborhood of the diagonal of $X^{q+1}$.
(14.16) Proposition. If $S$ is open or strongly paracompact in $X$, the relative cohomology groups can be written in terms of cohomology groups with supports in $X \backslash S$ :

$$
H^{q}(X, S ; M) \simeq H_{X \backslash S}^{q}(X, M)
$$

In particular, if $X \backslash S$ is relatively compact in $X$, we have

$$
H^{q}(X, S ; M) \simeq H_{c}^{q}(X \backslash S, M)
$$

Proof. We have an exact sequence

$$
\begin{equation*}
0 \longrightarrow M_{X \backslash S}^{[\bullet]}(X) \longrightarrow M^{[\bullet]}(X) \longrightarrow M^{[\bullet]}(S) \longrightarrow 0 \tag{14.17}
\end{equation*}
$$

where $M_{X \backslash S}^{[\bullet]}(X)$ denotes sections with support in $X \backslash S$. If $S$ is open, then $M^{[\bullet]}(S)=$ $\left(M_{\upharpoonright S}\right)^{[\bullet]}(S)$, hence $M_{X \backslash S}^{[\bullet]}(X)=M^{[\bullet]}(X, S)$ and the result follows. If $S$ is strongly paracompact, Prop. 4.7 and Th. 9.10 show that

$$
H^{q}\left(M^{[\bullet]}(S)\right)=H^{q}\left(\underset{\Omega \supset S}{\lim } M^{[\bullet]}(\Omega)\right)=\underset{\Omega \supset S}{\lim _{\Omega}} H^{q}(\Omega, M)=H^{q}\left(S, M_{\upharpoonright S}\right)
$$

If we consider the diagram

$$
\begin{array}{ccc}
0 \longrightarrow M_{X \backslash S}^{[\bullet]}(X) & \longrightarrow M^{[\bullet]}(X) & \longrightarrow M^{[\bullet]}(S)
\end{array} \longrightarrow 00
$$

we see that the last two vertical arrows induce isomorphisms in cohomology. Therefore, the first one also does.
(14.18) Corollary. Let $X, Y$ be locally compact spaces and $f, g: X \rightarrow Y$ proper maps. We say that $f, g$ are properly homotopic if they are homotopic through a proper homotopy $H: X \times I \longrightarrow Y$. Then

$$
f^{\star}=g^{\star}: \quad H_{c}^{q}(Y, M) \longrightarrow H_{c}^{q}(X, M) .
$$

Proof. Let $\widehat{X}=X \cup\{\infty\}, \widehat{Y}=Y \cup\{\infty\}$ be the Alexandrov compactifications of $X, Y$. Then $f, g, H$ can be extended as continuous maps

$$
\widehat{f}, \widehat{g}: \quad \widehat{X} \longrightarrow \widehat{Y}, \quad \widehat{H}: \quad \widehat{X} \times I \longrightarrow \widehat{Y}
$$

with $\widehat{f}(\infty)=\widehat{g}(\infty)=H(\infty, t)=\infty$, so that $\widehat{f}, \widehat{g}$ are homotopic as maps $(\widehat{X}, \infty) \longrightarrow$ $(\widehat{Y}, \infty)$. Proposition 14.16 implies $H_{c}^{q}(X, M)=H^{q}(\widetilde{X}, \infty ; M)$ and the result follows.

## § 15. Künneth Formula

## $\S$ 15.A. Flat Modules and Tor Functors

The goal of this section is to investigate homological properties related to tensor products. We work in the category of modules over a commutative ring $R$ with unit. All tensor products appearing here are tensor products over $R$. The starting point is the observation that tensor product with a given module is a right exact functor: if $0 \rightarrow A \rightarrow B \rightarrow C \rightarrow 0$ is an exact sequence and $M$ a $R$-module, then

$$
A \otimes M \longrightarrow B \otimes M \longrightarrow C \otimes M \longrightarrow 0
$$

is exact, but the map $A \otimes M \longrightarrow B \otimes M$ need not be injective. A counterexample is given by the sequence

$$
0 \longrightarrow \mathbb{Z} \xrightarrow{2 \times} \mathbb{Z} \longrightarrow \mathbb{Z} / 2 \mathbb{Z} \longrightarrow 0 \quad \text { over } \quad R=\mathbb{Z}
$$

tensorized by $M=\mathbb{Z} / 2 \mathbb{Z}$. However, the injectivity holds if $M$ is a free $R$-module. More generally, one says that $M$ is a flat $R$-module if the tensor product by $M$ preserves injectivity, or equivalently, if $\otimes M$ is a left exact functor.

A flat resolution $C \bullet$ of a $R$-module $A$ is a homology exact sequence

$$
\cdots \longrightarrow C_{q} \longrightarrow C_{q-1} \longrightarrow \cdots \longrightarrow C_{1} \longrightarrow C_{0} \longrightarrow A \longrightarrow 0
$$

where $C_{q}$ are flat $R$-modules and $C_{q}=0$ for $q<0$. Such a resolution always exists because every module $A$ is a quotient of a free module $C_{0}$. Inductively, we take $C_{q+1}$ to be a free module such that $\operatorname{ker}\left(C_{q} \rightarrow C_{q-1}\right)$ is a quotient of $C_{q+1}$. In terms of homology groups, we have $H_{0}\left(C_{\bullet}\right)=A$ and $H_{q}\left(C_{\bullet}\right)=0$ for $q \neq 0$. Given $R$-modules $A, B$ and free resolutions $d^{\prime}: C \bullet \longrightarrow A, d^{\prime \prime}: D \bullet \longrightarrow B$, we consider the double homology complex

$$
K_{p, q}=C_{p} \otimes D_{q}, \quad d=d^{\prime} \otimes \operatorname{Id}+(-1)^{p} \operatorname{Id} \otimes d^{\prime \prime}
$$

and the associated first and second spectral sequences. Since $C_{p}$ is free, we have

$$
E_{p, q}^{1}=H_{q}\left(C_{p} \otimes D_{\bullet}\right)= \begin{cases}C_{p} \otimes B & \text { for } q=0, \\ 0 & \text { for } q \neq 0\end{cases}
$$

Similarly, the second spectral sequence also collapses and we have

$$
H_{l}\left(K_{\bullet}\right)=H_{l}\left(C_{\bullet} \otimes B\right)=H_{l}\left(A \otimes D_{\bullet}\right) .
$$

This implies in particular that the homology groups $H_{l}\left(K_{\bullet}\right)$ do not depend on the choice of the resolutions $C_{\bullet}$ or $D_{\bullet}$.
(15.1) Definition. The $q$-th torsion module of $A$ and $B$ is

$$
\operatorname{Tor}_{q}(A, B)=H_{q}\left(K_{\bullet}\right)=H_{q}(C \bullet B)=H_{q}\left(A \otimes D_{\bullet}\right) .
$$

Since the definition of $K_{\bullet}$ is symmetric with respect to $A, B$, we have $\operatorname{Tor}_{q}(A, B) \simeq$ $\operatorname{Tor}_{q}(B, A)$. By the right-exactness of $\otimes B$, we find in particular $\operatorname{Tor}_{0}(A, B)=A \otimes B$. Moreover, if $B$ is flat, $\otimes B$ is also left exact, thus $\operatorname{Tor}_{q}(A, B)=0$ for all $q \geqslant 1$ and all modules $A$. If $0 \rightarrow A \rightarrow A^{\prime} \rightarrow A^{\prime \prime} \rightarrow 0$ is an exact sequence, there is a corresponding exact sequence of homology complexes

$$
0 \longrightarrow A \otimes D_{\bullet} \longrightarrow A^{\prime} \otimes D_{\bullet} \longrightarrow A^{\prime \prime} \otimes D_{\bullet} \longrightarrow 0
$$

thus a long exact sequence

$$
\begin{align*}
& \longrightarrow \operatorname{Tor}_{q}(A, B) \longrightarrow \operatorname{Tor}_{q}\left(A^{\prime}, B\right) \longrightarrow \operatorname{Tor}_{q}\left(A^{\prime \prime}, B\right) \longrightarrow \operatorname{Tor}_{q-1}(A, B) \\
& \cdots \longrightarrow A \otimes B \longrightarrow A^{\prime} \otimes B \longrightarrow A \tag{15.2}
\end{align*}
$$

It follows that $B$ is flat if and only if $\operatorname{Tor}_{1}(A, B)=0$ for every $R$-module $A$.
Suppose now that $R$ is a principal ring. Then every module $A$ has a free resolution $0 \rightarrow C_{1} \rightarrow C_{0} \rightarrow A \rightarrow 0$ because the kernel of any surjective map $C_{0} \rightarrow A$ is free (every submodule of a free module is free). It follows that one always has $\operatorname{Tor}_{q}(A, B)=0$ for $q \geqslant 2$. In this case, we denote $\operatorname{Tor}_{1}(A, B)=A \star B$ and call it the torsion product of $A$ and $B$. The above exact sequence ( 15.2 ) reduces to

$$
\begin{equation*}
0 \rightarrow A \star B \rightarrow A^{\prime} \star B \rightarrow A^{\prime \prime} \star B \rightarrow A \otimes B \rightarrow A^{\prime} \otimes B \rightarrow A^{\prime \prime} \otimes B \rightarrow 0 \tag{15.3}
\end{equation*}
$$

In order to compute $A \star B$, we may restrict ourselves to finitely generated modules, because every module is a direct limit of such modules and the $\star$ product commutes with direct limits. Over a principal ring $R$, every finitely generated module is a direct sum of a free module and of cyclic modules $R / a R$. It is thus sufficient to compute $R / a R \star R / b R$. The obvious free resolution $R \xrightarrow{a \times} R$ of $R / a R$ shows that $R / a R \star R / b R$ is the kernel of the map $R / b R \xrightarrow{a \times} R / b R$. Hence

$$
\begin{equation*}
R / a R \star R / b R \simeq R /(a \wedge b) R \tag{15.4}
\end{equation*}
$$

where $a \wedge b$ denotes the greatest common divisor of $a$ and $b$. It follows that a module $B$ is flat if and only if it is torsion free. If $R$ is a field, every $R$-module is free, thus $A \star B=0$ for all $A$ and $B$.

## §15.B. Künneth and Universal Coefficient Formulas

The algebraic Künneth formula describes the cohomology groups of the tensor product of two differential complexes.
(15.5) Algebraic Künneth formula. Let $\left(K^{\bullet}, d^{\prime}\right),\left(L^{\bullet}, d^{\prime \prime}\right)$ be complexes of $R$-modules and $(K \otimes L)^{\bullet}$ the simple complex associated to the double complex $(K \otimes L)^{p, q}=K^{p} \otimes L^{q}$. If $K^{\bullet}$ or $L^{\bullet}$ is torsion free, there is a split exact sequence

$$
0 \rightarrow \bigoplus_{p+q=l} H^{p}\left(K^{\bullet}\right) \otimes H^{q}\left(L^{\bullet}\right) \xrightarrow{\mu} H^{l}\left((K \otimes L)^{\bullet}\right) \rightarrow \bigoplus_{p+q=l+1} H^{p}\left(K^{\bullet}\right) \star H^{q}\left(L^{\bullet}\right)
$$

where the map $\mu$ is defined by $\mu\left(\left\{k^{p}\right\} \times\left\{l^{q}\right\}\right)=\left\{k^{p} \otimes l^{q}\right\}$ for all cocycles $\left\{k^{p}\right\} \in H^{p}\left(K^{\bullet}\right)$, $\left\{l^{q}\right\} \in H^{q}\left(L^{\bullet}\right)$.
(15.6) Corollary. If $R$ is a field, or if one of the graded modules $H^{\bullet}\left(K^{\bullet}\right), H^{\bullet}\left(L^{\bullet}\right)$ is torsion free, then

$$
H^{l}\left((K \otimes L)^{\bullet}\right) \simeq \bigoplus_{p+q=l} H^{p}\left(K^{\bullet}\right) \otimes H^{q}\left(L^{\bullet}\right)
$$

Proof. Assume for example that $K^{\bullet}$ is torsion free. There is a short exact sequence of complexes

$$
0 \longrightarrow Z^{\bullet} \longrightarrow K^{\bullet} \xrightarrow{d^{\prime}} B^{\bullet+1} \longrightarrow 0
$$

where $Z^{\bullet}, B^{\bullet} \subset K^{\bullet}$ are respectively the graded modules of cocycles and coboundaries in $K^{\bullet}$, considered as subcomplexes with zero differential. As $B^{\bullet+1}$ is torsion free, the tensor product of the above sequence with $L^{\bullet}$ is still exact. The corresponding long exact sequence for the associated simple complexes yields:

$$
\begin{align*}
H^{l}\left((B \otimes L)^{\bullet}\right) & \longrightarrow H^{l}\left((Z \otimes L)^{\bullet}\right) \longrightarrow H^{l}\left((K \otimes L)^{\bullet}\right) \xrightarrow{d^{\prime}} H^{l+1}\left((B \otimes L)^{\bullet}\right) \\
& \longrightarrow H^{l+1}\left((Z \otimes L)^{\bullet}\right) \cdots \tag{15.7}
\end{align*}
$$

The first and last arrows are connecting homomorphisms; in this situation, they are easily seen to be induced by the inclusion $B^{\bullet} \subset Z^{\bullet}$. Since the differential of $Z^{\bullet}$ is zero, the simple complex $(Z \otimes L)^{\bullet}$ is isomorphic to the direct sum $\bigoplus_{p} Z^{p} \otimes L^{\bullet-p}$, where $Z^{p}$ is torsion free. Similar properties hold for $(B \otimes L)^{\bullet}$, hence

$$
H^{l}\left((Z \otimes L)^{\bullet}\right)=\bigoplus_{p+q=l} Z^{p} \otimes H^{q}\left(L^{\bullet}\right), \quad H^{l}\left((B \otimes L)^{\bullet}\right)=\bigoplus_{p+q=l} B^{p} \otimes H^{q}\left(L^{\bullet}\right)
$$

The exact sequence

$$
0 \longrightarrow B^{p} \longrightarrow Z^{p} \longrightarrow H^{p}\left(K^{\bullet}\right) \longrightarrow 0
$$

tensorized by $H^{q}\left(L^{\bullet}\right)$ yields an exact sequence of the type (15.3):

$$
\begin{aligned}
0 \rightarrow H^{p}\left(K^{\bullet}\right) \star H^{q}\left(L^{\bullet}\right) \rightarrow B^{p} \otimes & H^{q}\left(L^{\bullet}\right) \rightarrow Z^{p} \otimes H^{q}\left(L^{\bullet}\right) \\
& \rightarrow H^{p}\left(K^{\bullet}\right) \otimes H^{q}\left(L^{\bullet}\right) \rightarrow 0 .
\end{aligned}
$$

By the above equalities, we get

$$
\begin{array}{r}
0 \longrightarrow \bigoplus_{p+q=l} H^{p}\left(K^{\bullet}\right) \star H^{q}\left(L^{\bullet}\right) \longrightarrow H^{l}\left((B \otimes L)^{\bullet}\right) \longrightarrow H^{l}\left((Z \otimes L)^{\bullet}\right) \\
\longrightarrow \bigoplus_{p+q=l} H^{p}\left(K^{\bullet}\right) \otimes H^{q}\left(L^{\bullet}\right) \longrightarrow 0 .
\end{array}
$$

In our initial long exact sequence (15.7), the cokernel of the first arrow is thus

$$
\bigoplus_{p+q=l} H^{p}\left(K^{\bullet}\right) \otimes H^{q}\left(L^{\bullet}\right)
$$

and the kernel of the last arrow is the torsion sum $\bigoplus_{p+q=l+1} H^{p}\left(K^{\bullet}\right) \star H^{q}\left(L^{\bullet}\right)$. This gives the exact sequence of the lemma. We leave the computation of the map $\mu$ as an exercise for the reader. The splitting assertion can be obtained by observing that there always exists a torsion free complex $\widetilde{K}^{\bullet}$ that splits (i.e. $\widetilde{Z}^{\bullet} \subset \widetilde{K}^{\bullet}$ splits), and a morphism $\widetilde{K}^{\bullet} \longrightarrow K^{\bullet}$ inducing an isomorphism in cohomology; then the projection $\widetilde{K}^{\bullet} \longrightarrow \widetilde{Z}^{\bullet}$ yields a projection

$$
\begin{aligned}
H^{l}\left((\widetilde{K} \otimes L)^{\bullet}\right) \longrightarrow H^{l}\left((\widetilde{Z} \otimes L)^{\bullet}\right) \simeq \bigoplus_{p+q=l} \widetilde{Z}^{p} & \otimes H^{q}\left(L^{\bullet}\right) \\
& \longrightarrow \bigoplus_{p+q=l} H^{p}\left(\widetilde{K}^{\bullet}\right) \otimes H^{q}\left(L^{\bullet}\right)
\end{aligned}
$$

To construct $\widetilde{K}^{\bullet}$, let $\widetilde{Z}^{\bullet} \longrightarrow Z^{\bullet}$ be a surjective map with $\widetilde{Z}^{\bullet}$ free, $\widetilde{B}^{\bullet}$ the inverse image of $B^{\bullet}$ in $\widetilde{\mathscr{C}}^{\bullet}$ and $\widetilde{K}^{\bullet}=\widetilde{Z}^{\bullet} \oplus \widetilde{B}^{\bullet+1}$, where the differential $\widetilde{K}^{\bullet} \longrightarrow \widetilde{K}^{\bullet+1}$ is given by $\widetilde{Z}^{\bullet} \longrightarrow 0$ and $\widetilde{B}^{\bullet+1} \subset \widetilde{Z}^{\bullet+1} \oplus 0 ;$ as $\widetilde{B}^{\bullet}$ is free, the map $\widetilde{B}^{\bullet+1} \longrightarrow B^{\bullet+1}$ can be lifted to a map $\widetilde{B}^{\bullet+1} \longrightarrow K^{\bullet}$, and this lifting combined with the composite $\widetilde{Z}^{\bullet} \rightarrow Z^{\bullet} \subset K^{\bullet}$ yields the required complex morphism $\widetilde{K}^{\bullet}=\widetilde{Z}^{\bullet} \oplus \widetilde{B}^{\bullet+1} \longrightarrow K^{\bullet}$.
(15.8) Universal coefficient formula. Let $K^{\bullet}$ be a complex of $R$-modules and $M a$ $R$-module such that either $K^{\bullet}$ or $M$ is torsion free. Then there is a split exact sequence

$$
0 \longrightarrow H^{p}\left(K^{\bullet}\right) \otimes M \longrightarrow H^{p}\left(K^{\bullet} \otimes M\right) \longrightarrow H^{p+1}\left(K^{\bullet}\right) \star M \longrightarrow 0 .
$$

Indeed, this is a special case of Formula 15.5 when the complex $L^{\bullet}$ is reduced to one term $L^{0}=M$. In general, it is interesting to observe that the spectral sequence of $K^{\bullet} \otimes L^{\bullet}$ collapses in $E_{2}$ if $K^{\bullet}$ is torsion free: $H^{k}\left((K \otimes L)^{\bullet}\right)$ is in fact the direct sum of the terms $E_{2}^{p, q}=H^{p}\left(K^{\bullet} \otimes H^{q}\left(L^{\bullet}\right)\right)$ thanks to (15.8).

## § 15.C. Künneth Formula for Sheaf Cohomology

Here we apply the general algebraic machinery to compute cohomology groups over a product space $X \times Y$. The main argument is a combination of the Leray spectral sequence with the universal coefficient formula for sheaf cohomology.
(15.9) Theorem. Let $\mathscr{A}$ be a sheaf of $R$-modules over a topological space $X$ and $M a$ $R$-module. Assume that either $\mathscr{A}$ or $M$ is torsion free and that either $X$ is compact or $M$ is finitely generated. Then there is a split exact sequence

$$
0 \longrightarrow H^{p}(X, \mathscr{A}) \otimes M \longrightarrow H^{p}(X, \mathscr{A} \otimes M) \longrightarrow H^{p+1}(X, \mathscr{A}) \star M \longrightarrow 0 .
$$

Proof. If $M$ is finitely generated, we get $(\mathscr{A} \otimes M)^{[\bullet]}(X)=\mathscr{A}^{[\bullet]}(X) \otimes M$ easily, so the above exact sequence is a consequence of Formula 15.8. If $X$ is compact, we may consider Čech cochains $C^{q}(U, \mathscr{A} \otimes M)$ over finite coverings. There is an obvious morphism

$$
C^{q}(\mathscr{U}, \mathscr{A}) \otimes M \longrightarrow C^{q}(\mathscr{U}, \mathscr{A} \otimes M)
$$

but this morphism need not be surjective nor injective. However, since

$$
(\mathscr{A} \otimes M)_{x}=\mathscr{A}_{x} \otimes M=\underset{V \ni x}{\lim } \mathscr{A}(V) \otimes M,
$$

the following properties are easy to verify:
a) If $c \in C^{q}(\mathscr{U}, \mathscr{A} \otimes M)$, there is a refinement $\mathscr{V}$ of $\mathscr{U}$ and $\rho: \mathscr{V} \longrightarrow \mathscr{U}$ such that $\rho^{\star} c \in C^{q}(\mathscr{V}, \mathscr{A} \otimes M)$ is in the image of $C^{q}(\mathscr{V}, \mathscr{A}) \otimes M$.
b) If a tensor $t \in C^{q}(\mathscr{U}, \mathscr{A}) \otimes M$ is mapped to 0 in $C^{q}(\mathscr{U}, \mathscr{A} \otimes M)$, there is a refinement $\mathscr{V}$ of $\mathscr{U}$ such that $\rho^{\star} t \in C^{q}(\mathscr{V}, \mathscr{A}) \otimes M$ equals 0 .

From a) and b) it follows that

$$
\check{H}^{q}(X, \mathscr{A} \otimes M)=\underset{\ddots}{\lim } H^{q}\left(C^{\bullet}(\mathscr{U}, \mathscr{A} \otimes M)\right)=\underset{\vartheta}{\lim } H^{q}\left(C^{\bullet}(\mathscr{U}, \mathscr{A}) \otimes M\right)
$$

and the desired exact sequence is the direct limit of the exact sequences of Formula 15.8 obtained for $K^{\bullet}=C^{\bullet}(\mathscr{U}, \mathscr{A})$.
(15.10) Theorem (Künneth). Let $\mathscr{A}$ and $\mathscr{B}$ be sheaves of $R$-modules over topological spaces $X$ and $Y$. Assume that $\mathscr{A}$ is torsion free, that $Y$ is compact and that either $X$ is compact or the cohomology groups $H^{q}(Y, \mathscr{B})$ are finitely generated $R$-modules. There is a split exact sequence

$$
\begin{aligned}
0 \longrightarrow \bigoplus_{p+q=l} H^{p}(X, \mathscr{A}) & \otimes H^{q}(Y, \mathscr{B}) \xrightarrow{\mu} H^{l}(X \times Y, \mathscr{A} \boxtimes \mathscr{B}) \\
& \longrightarrow \bigoplus_{p+q=l+1} H^{p}(X, \mathscr{A}) \star H^{q}(Y, \mathscr{B}) \longrightarrow 0
\end{aligned}
$$

where $\mu$ is the map given by the cartesian product $\bigoplus \alpha_{p} \otimes \beta_{q} \longmapsto \sum \alpha_{p} \times \beta_{q}$.
Proof. We compute $H^{l}(X, \mathscr{A} \boxtimes \mathscr{B})$ by means of the Leray spectral sequence of the projection $\pi: X \times Y \longrightarrow X$. This means that we are considering the differential sheaf $\mathscr{L}^{q}=\pi_{\star}(\mathscr{A} \boxtimes \mathscr{B})^{[q]}$ and the double complex

$$
K^{p, q}=\left(\mathscr{L}^{q}\right)^{[p]}(X)
$$

By (12.5') we have ${ }_{K} E_{2}^{p, q}=H^{p}\left(X, \mathscr{H}^{q}\left(\mathscr{L}^{\bullet}\right)\right)$. As $Y$ is compact, the cohomology sheaves $\mathscr{H}^{q}\left(\mathscr{L}^{\bullet}\right)=R^{q} \pi_{\star}(\mathscr{A} \boxtimes \mathscr{B})$ are given by

$$
R^{q} \pi_{\star}(\mathscr{A} \boxtimes \mathscr{B})_{x}=H^{q}\left(\{x\} \times Y, \mathscr{A} \boxtimes \mathscr{B}_{\mid\{x\} \times Y}\right)=H^{q}\left(Y, \mathscr{A}_{x} \otimes \mathscr{B}\right)=\mathscr{A}_{x} \otimes H^{q}(Y, \mathscr{B})
$$

thanks to the compact case of Th. 15.9 where $M=\mathscr{A}_{x}$ is torsion free. We obtain therefore

$$
\begin{aligned}
& R^{q} \pi_{\star}(\mathscr{A} \boxtimes \mathscr{B})=\mathscr{A} \otimes H^{q}(Y, \mathscr{B}), \\
& { }_{K} E_{2}^{p, q}=H^{p}\left(X, \mathscr{A} \otimes H^{q}(Y, \mathscr{B})\right) .
\end{aligned}
$$

Theorem 15.9 shows that the $E_{2}$-term is actually given by the desired exact sequence, but it is not a priori clear that the spectral sequence collapses in $E_{2}$. In order to check this, we consider the double complex

$$
C^{p, q}=\mathscr{A}^{[p]}(X) \otimes \mathscr{B}^{[q]}(Y)
$$

and construct a natural morphism $C^{\bullet \bullet \bullet} \longrightarrow K^{\bullet \bullet \bullet}$. We may consider the set $K^{p, q}=$ $\left(\pi_{\star}(\mathscr{A} \boxtimes \mathscr{B})^{[q]}\right)^{[p]}(X)$ as the set of equivalence classes of functions

$$
h\left(\xi_{0}, \ldots, \xi_{p}\right) \in \pi_{\star}(\mathscr{A} \boxtimes \mathscr{B})_{\xi_{p}}^{[q]}=\underset{\longrightarrow}{\lim }(\mathscr{A} \boxtimes \mathscr{B})^{[q]}\left(\pi^{-1}\left(V\left(\xi_{p}\right)\right)\right)
$$

or more precisely

$$
\begin{aligned}
& h\left(\xi_{0}, \ldots, \xi_{p} ;\left(x_{0}, y_{0}\right), \ldots,\left(x_{q}, y_{q}\right)\right) \in \mathscr{A}_{x_{q}} \otimes \mathscr{B}_{y_{q}} \quad \text { with } \\
& \xi_{0} \in X, \quad \xi_{j} \in V\left(\xi_{0}, \ldots, \xi_{j-1}\right), \quad 1 \leqslant j \leqslant p, \\
& \left(x_{0}, y_{0}\right) \in V\left(\xi_{0}, \ldots, \xi_{p}\right) \times Y, \\
& \left(x_{j}, y_{j}\right) \in V\left(\xi_{0}, \ldots, \xi_{p} ;\left(x_{0}, y_{0}\right), \ldots,\left(x_{j-1}, y_{j-1}\right)\right), \quad 1 \leqslant j \leqslant q .
\end{aligned}
$$

Then $f \otimes g \in C^{p, q}$ is mapped to $h \in K^{p, q}$ by the formula

$$
h\left(\xi_{0}, \ldots, \xi_{p} ;\left(x_{0}, y_{0}\right), \ldots,\left(x_{q}, y_{q}\right)\right)=f\left(\xi_{0}, \ldots, \xi_{p}\right)\left(x_{q}\right) \otimes g\left(y_{0}, \ldots, y_{q}\right)
$$

As $\mathscr{A}^{[p]}(X)$ is torsion free, we find

$$
{ }_{C} E_{1}^{p, q}=\mathscr{A}^{[p]}(X) \otimes H^{q}(Y, \mathscr{B})
$$

Since either $X$ is compact or $H^{q}(Y, \mathscr{B})$ finitely generated, Th. 15.9 yields

$$
{ }_{C} E_{2}^{p, q}=H^{p}\left(X, \mathscr{A} \otimes H^{q}(Y, \mathscr{B})\right) \simeq{ }_{K} E_{2}^{p, q}
$$

hence $H^{l}\left(K^{\bullet}\right) \simeq H^{l}\left(C^{\bullet}\right)$ and the algebraic Künneth formula 15.5 concludes the proof.
(15.11) Remark. The exact sequences of Th. 15.9 and of Künneth's theorem also hold for cohomology groups with compact support, provided that $X$ and $Y$ are locally compact and $\mathscr{A}$ (or $\mathscr{B}$ ) is torsion free. This is an immediate consequence of Prop. 7.12 on direct limits of cohomology groups over compact subsets.
(15.12) Corollary. When $\mathscr{A}$ and $\mathscr{B}$ are torsion free constant sheaves, e.g. $\mathscr{A}=\mathscr{B}=\mathbb{Z}$ or $\mathbb{R}$, the Künneth formula holds as soon as $X$ or $Y$ has the same homotopy type as a finite cell complex.

Proof. If $Y$ satisfies the assumption, we may suppose in fact that $Y$ is a finite cell complex by the homotopy invariance. Then $Y$ is compact and $H^{\bullet}(Y, \mathscr{B})$ is finitely generated, so Th. 15.10 can be applied.

## § 16. Poincaré duality

## §16.A. Injective Modules and Ext Functors

The study of duality requires some algebraic preliminaries on the Hom functor and its derived functors Ext ${ }^{q}$. Let $R$ be a commutative ring with unit, $M$ a $R$-module and

$$
0 \longrightarrow A \longrightarrow B \longrightarrow C \longrightarrow 0
$$

an exact sequence of $R$-modules. Then we have exact sequences

$$
\begin{aligned}
0 \longrightarrow & \operatorname{Hom}_{R}(M, A) \longrightarrow \operatorname{Hom}_{R}(M, B) \longrightarrow \operatorname{Hom}_{R}(M, C), \\
& \operatorname{Hom}_{R}(A, M) \longleftarrow \operatorname{Hom}_{R}(B, M) \longleftarrow \operatorname{Hom}_{R}(C, M) \longleftarrow 0
\end{aligned}
$$

i.e. $\operatorname{Hom}(M, \bullet)$ is a covariant left exact functor and $\operatorname{Hom}(\bullet, M)$ a contravariant right exact functor. The module $M$ is said to be projective if $\operatorname{Hom}(M, \bullet)$ is also right exact, and injective if $\operatorname{Hom}(\bullet, M)$ is also left exact. Every free $R$-module is projective. Conversely, if $M$ is projective, any surjective morphism $F \longrightarrow M$ from a free module $F$ onto $M$ must split $\left(\operatorname{Id}_{M}\right.$ has a preimage in $\left.\operatorname{Hom}(M, F)\right)$; if $R$ is a principal ring, "projective" is therefore equivalent to "free".
(16.1) Proposition. Over a principal ring $R$, a module $M$ is injective if and only if it is divisible, i.e. if for every $x \in M$ and $\lambda \in R \backslash\{0\}$, there exists $y \in M$ such that $\lambda y=x$.

Proof. If $M$ is injective, the exact sequence $0 \longrightarrow R \xrightarrow{\lambda \times} R \longrightarrow R / \lambda R \longrightarrow 0$ shows that

$$
M=\operatorname{Hom}(R, M) \xrightarrow{\lambda \times} \operatorname{Hom}(R, M)=M
$$

must be surjective, thus $M$ is divisible.
Conversely, assume that $R$ is divisible. Let $f: A \longrightarrow M$ be a morphism and $B \supset A$. Zorn's lemma implies that there is a maximal extension $\widetilde{\sim}: \widetilde{A} \longrightarrow M$ of $f$ where $A \subset$ $\widetilde{A} \subset B \dot{\tilde{A}}$. If $\widetilde{A} \neq B$, select $x \in B \backslash \widetilde{A}$ and consider the ideal $I$ of elements $\lambda \in R$ such that $\lambda x \in \widetilde{A}$. As $R$ is principal we have $I=\lambda_{0} R$ for some $\lambda_{0}$. If $\lambda_{0} \neq 0$, select $y \in M$ such that $\lambda_{0} y=\widetilde{f}\left(\lambda_{0} x\right)$ and if $\lambda_{0}=0$ take $y$ arbitrary. Then $\widetilde{f}$ can be extended to $\widetilde{A}+R x$ by letting $\widetilde{f}(x)=y$. This is a contradiction, so we must have $\widetilde{A}=B$.
(16.2) Proposition. Every module $M$ can be embedded in an injective module $\widetilde{M}$.

Proof. Assume first $R=\mathbb{Z}$. Then set

$$
M^{\prime}=\operatorname{Hom}_{\mathbb{Z}}(M, \mathbb{Q} / \mathbb{Z}), \quad M^{\prime \prime}=\operatorname{Hom}_{\mathbb{Z}}\left(M^{\prime}, \mathbb{Q} / \mathbb{Z}\right) \subset \mathbb{Q} / \mathbb{Z}^{M^{\prime}}
$$

Since $\mathbb{Q} / \mathbb{Z}$ is divisible, $\mathbb{Q} / \mathbb{Z}$ and $\mathbb{Q} / \mathbb{Z}^{M^{\prime}}$ are injective. It is therefore sufficient to show that the canonical morphism $M \longrightarrow M^{\prime \prime}$ is injective. In fact, for any $x \in M \backslash\{0\}$, the subgroup $\mathbb{Z} x$ is cyclic (finite or infinite), so there is a non trivial morphism $\mathbb{Z} x \longrightarrow \mathbb{Q} / \mathbb{Z}$, and we can extend this morphism into a morphism $u: M \longrightarrow \mathbb{Q} / \mathbb{Z}$. Then $u \in M^{\prime}$ and $u(x) \neq 0$, so $M \longrightarrow M^{\prime \prime}$ is injective.

Now, for an arbitrary ring $R$, we set $\widetilde{M}=\operatorname{Hom}_{\mathbb{Z}}\left(R, \mathbb{Q} / \mathbb{Z}^{M^{\prime}}\right)$. There are $R$-linear embeddings

$$
M=\operatorname{Hom}_{R}(R, M) \hookrightarrow \operatorname{Hom}_{\mathbb{Z}}(R, M) \hookrightarrow \operatorname{Hom}_{\mathbb{Z}}\left(R, \mathbb{Q} / \mathbb{Z}^{M^{\prime}}\right)=\widetilde{M}
$$

and since $\operatorname{Hom}_{R}(\bullet, \widetilde{M}) \simeq \operatorname{Hom}_{\mathbb{Z}}\left(\bullet, \mathbb{Q} / \mathbb{Z}^{M^{\prime}}\right)$, it is clear that $\widetilde{M}$ is injective over the ring $R$.

As a consequence, any module has a (cohomological) resolution by injective modules. Let $A, B$ be given $R$-modules, let $d^{\prime}: B \rightarrow D^{\bullet}$ be an injective resolution of $B$ and let $d^{\prime \prime}: C \bullet \rightarrow A$ be a free (or projective) resolution of $A$. We consider the cohomology double complex

$$
K^{p, q}=\operatorname{Hom}\left(C_{q}, D^{p}\right), \quad d=d^{\prime}+(-1)^{p}\left(d^{\prime \prime}\right)^{\dagger}
$$

( $\dagger$ means transposition) and the associated first and second spectral sequences. Since $\operatorname{Hom}\left(\bullet, D^{p}\right)$ and $\operatorname{Hom}\left(C_{q}, \bullet\right)$ are exact, we get

$$
\begin{aligned}
& E_{1}^{p, 0}=\operatorname{Hom}\left(A, D^{p}\right), \quad \widetilde{E}_{1}^{p, 0}=\operatorname{Hom}\left(C_{p}, B\right), \\
& E_{1}^{p, q}=\widetilde{E}_{1}^{p, q}=0 \quad \text { for } \quad q \neq 0 .
\end{aligned}
$$

Therefore, both spectral sequences collapse in $E_{1}$ and we get

$$
H^{l}\left(K^{\bullet}\right)=H^{l}\left(\operatorname{Hom}\left(A, D^{\bullet}\right)\right)=H^{l}(\operatorname{Hom}(C \bullet, B)) ;
$$

in particular, the cohomology groups $H^{l}\left(K^{\bullet}\right)$ do not depend on the choice of the resolutions $C$ • or $D^{\bullet}$.
(16.3) Definition. The $q$-th extension module of $A, B$ is

$$
\operatorname{Ext}_{R}^{q}(A, B)=H^{q}\left(K^{\bullet}\right)=H^{q}\left(\operatorname{Hom}\left(A, D^{\bullet}\right)\right)=H^{q}\left(\operatorname{Hom}\left(C_{\bullet}, B\right)\right)
$$

By the left exactness of $\operatorname{Hom}(A, \bullet)$, we get in particular $\operatorname{Ext}^{0}(A, B)=\operatorname{Hom}(A, B)$. If $A$ is projective or $B$ injective, then clearly $\operatorname{Ext}^{q}(A, B)=0$ for all $q \geqslant 1$. Any exact sequence $0 \rightarrow A \rightarrow A^{\prime} \rightarrow A^{\prime \prime} \rightarrow 0$ is converted into an exact sequence by $\operatorname{Hom}\left(\bullet, D^{\bullet}\right)$, thus we get a long exact sequence

$$
\begin{aligned}
0 & \longrightarrow \operatorname{Hom}\left(A^{\prime \prime}, B\right) \longrightarrow \operatorname{Hom}\left(A^{\prime}, B\right) \longrightarrow \operatorname{Hom}(A, B) \longrightarrow \operatorname{Ext}^{1}\left(A^{\prime \prime}, B\right) \cdots \\
& \longrightarrow \operatorname{Ext}^{q}\left(A^{\prime \prime}, B\right) \longrightarrow \operatorname{Ext}^{q}\left(A^{\prime}, B\right) \longrightarrow \operatorname{Ext}^{q}(A, B) \longrightarrow \operatorname{Ext}^{q+1}\left(A^{\prime \prime}, B\right) \cdots
\end{aligned}
$$

Similarly, any exact sequence $0 \rightarrow B \rightarrow B^{\prime} \rightarrow B^{\prime \prime} \rightarrow 0$ yields

$$
\begin{aligned}
0 & \longrightarrow \operatorname{Hom}(A, B) \longrightarrow \operatorname{Hom}\left(A, B^{\prime}\right) \longrightarrow \operatorname{Hom}\left(A, B^{\prime \prime}\right) \longrightarrow \operatorname{Ext}^{1}(A, B) \cdots \\
& \longrightarrow \operatorname{Ext}^{q}(A, B) \longrightarrow \operatorname{Ext}^{q}\left(A, B^{\prime}\right) \longrightarrow \operatorname{Ext}^{q}\left(A, B^{\prime \prime}\right) \longrightarrow \operatorname{Ext}^{q+1}(A, B) \cdots
\end{aligned}
$$

Suppose now that $R$ is a principal ring. Then the resolutions $C_{\bullet}$ or $D^{\bullet}$ can be taken of length 1 (any quotient of a divisible module is divisible), thus $\operatorname{Ext}^{q}(A, B)$ is always 0 for $q \geqslant 2$. In this case, we simply denote $\operatorname{Ext}^{1}(A, B)=\operatorname{Ext}(A, B)$. When $A$ is finitely generated, the computation of $\operatorname{Ext}(A, B)$ can be reduced to the cyclic case, since $\operatorname{Ext}(A, B)=0$ when $A$ is free. For $A=R / a R$, the obvious free resolution $R \xrightarrow{a \times} R$ gives

$$
\begin{equation*}
\operatorname{Ext}_{R}(R / a R, B)=B / a B \tag{16.4}
\end{equation*}
$$

(16.5) Lemma. Let $K_{\bullet}$ be a homology complex and let $M \rightarrow M^{\bullet}$ be an injective resolution of a $R$-module $M$. Let $L^{\bullet}$ be the simple complex associated to $L^{p, q}=\operatorname{Hom}_{R}\left(K_{q}, M^{p}\right)$. There is a split exact sequence

$$
0 \longrightarrow \operatorname{Ext}\left(H_{q-1}\left(K_{\bullet}\right), M\right) \longrightarrow H^{q}\left(L^{\bullet}\right) \longrightarrow \operatorname{Hom}\left(H_{q}\left(K_{\bullet}\right), M\right) \longrightarrow 0
$$

Proof. As the functor $\operatorname{Hom}_{R}\left(\bullet, M^{p}\right)$ is exact, we get

$$
\begin{aligned}
& { }_{L} E_{1}^{p, q}=\operatorname{Hom}\left(H_{q}\left(K_{\bullet}\right), M^{p}\right), \\
& { }_{L} E_{2}^{p, q}= \begin{cases}\operatorname{Hom}\left(H_{q}\left(K_{\bullet}\right), M\right) & \text { for } p=0, \\
\operatorname{Ext}\left(H_{q}\left(K_{\bullet}\right), M\right) & \text { for } p=1, \\
0 & \text { for } p \geqslant 2 .\end{cases}
\end{aligned}
$$

The spectral sequence collapses in $E_{2}$, therefore we get

$$
\begin{aligned}
& G_{0}\left(H^{q}\left(L^{\bullet}\right)\right)=\operatorname{Hom}\left(H_{q}\left(K_{\bullet}\right), M\right), \\
& G_{1}\left(H^{q}\left(L^{\bullet}\right)\right)=\operatorname{Ext}\left(H_{q-1}\left(K_{\bullet}\right), M\right)
\end{aligned}
$$

and the expected exact sequence follows. By the same arguments as at the end of the proof of Formula 15.5, we may assume that $K_{\bullet}$ is split, so that there is a projection $K_{q} \longrightarrow Z_{q}$. Then the composite morphism

$$
\begin{aligned}
\operatorname{Hom}\left(H_{q}\left(K_{\bullet}\right), M\right)=\operatorname{Hom}\left(Z_{q} / B_{q}, M\right) \longrightarrow \operatorname{Hom}( & \left.K_{q} / B_{q}, M\right) \\
& \subset Z^{q}\left(L^{\bullet}\right) \longrightarrow H^{q}\left(L^{\bullet}\right)
\end{aligned}
$$

defines a splitting of the exact sequence.

## § 16.B. Poincaré Duality for Sheaves

Let $\mathscr{A}$ be a sheaf of abelian groups on a locally compact topological space $X$ of finite topological dimension $n=\operatorname{topdim} X$. By 13.12 c ), $\mathscr{A}$ admits a soft resolution $\mathscr{L} \bullet$ of length $n$. If $M \rightarrow M^{0} \rightarrow M^{1} \rightarrow 0$ is an injective resolution of a $R$-module $M$, we introduce the double complex of presheaves $\mathscr{F}_{\mathscr{A}, M}^{p, q}$ defined by

$$
\begin{equation*}
\mathscr{F}_{\mathscr{A}, M}^{p, q}(U)=\operatorname{Hom}_{R}\left(\mathscr{L}_{c}^{n-q}(U), M^{p}\right), \tag{16.6}
\end{equation*}
$$

where the restriction map $\mathscr{F}_{A, M}^{p, q}(U) \longrightarrow \mathscr{F}_{\mathscr{A}, M}^{p, q}(V)$ is the adjoint of the inclusion map $\mathscr{L}_{c}^{n-q}(V) \longrightarrow \mathscr{L}_{c}^{n-q}(U)$ when $V \subset U$. As $\mathscr{L}^{n-q}$ is soft, any $f \in \mathscr{L}_{c}^{n-q}(U)$ can be written as $f=\sum f_{\alpha}$ with $\left(f_{\alpha}\right)$ subordinate to any open covering $\left(U_{\alpha}\right)$ of $U$; it follows easily that $\mathscr{F}_{\mathscr{A}, M}^{p, q}$ satisfy axioms (II-2.4) of sheaves. The injectivity of $M^{p}$ implies that $\mathscr{F}_{\mathscr{A}, M}^{p, q}$ is a flabby sheaf. By Lemma 16.5, we get a split exact sequence

$$
\begin{align*}
0 \longrightarrow \operatorname{Ext}\left(H_{c}^{n-q+1}(X, \mathscr{A}), M\right) \longrightarrow H^{q} & \left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}(X)\right) \\
& \longrightarrow \operatorname{Hom}\left(H_{c}^{n-q}(X, \mathscr{A}), M\right) \longrightarrow 0 . \tag{16.7}
\end{align*}
$$

This can be seen as an abstract Poincaré duality formula, relating the cohomology groups of a differential sheaf $\mathscr{F} \mathscr{A}, M$ "dual" of $\mathscr{A}$ to the dual of the cohomology with compact support of $\mathscr{A}$. In concrete applications, it still remains to compute $H^{q}\left(\mathscr{F} \bullet \mathscr{A}_{, M}(X)\right)$. This can be done easily when $X$ is a manifold and $\mathscr{A}$ is a constant or locally constant sheaf.

## § 16.C. Poincaré Duality on Topological Manifolds

Here, $X$ denotes a paracompact topological manifold of dimension $n$.
(16.8) Definition. Let $L$ be a $R$-module. A locally constant sheaf of stalk $L$ on $X$ is a sheaf $\mathscr{A}$ such that every point has a neighborhood $\Omega$ on which $\mathscr{A}_{\Gamma}$ is $R$-isomorphic to the constant sheaf $L$.

Thus, a locally constant sheaf $\mathscr{A}$ can be seen as a discrete fiber bundle over $X$ whose fibers are $R$-modules and whose transition automorphisms are $R$-linear. If $X$ is locally contractible, a locally constant sheaf of stalk $L$ is given, up to isomorphism, by a representation $\rho: \pi_{1}(X) \longrightarrow \operatorname{Aut}_{R}(L)$ of the fundamental group of $X$, up to conjugation; denoting by $\widetilde{X}$ the universal covering of $X$, the sheaf $\mathscr{A}$ associated to $\rho$ can be viewed as the quotient of $\widetilde{X} \times L$ by the diagonal action of $\pi_{1}(X)$. We leave the reader check himself the details of these assertions: in fact similar arguments will be explained in full details in §V-6 when properties of flat vector bundles are discussed.

Let $\mathscr{A}$ be a locally constant sheaf of stalk $L$, let $\mathscr{L} \bullet$ be a soft resolution of $\mathscr{A}$ and $\mathscr{F}_{A, M}^{p, q}$ the associated flabby sheaves. For an arbitrary open set $U \subset X$, Formula (16.7) gives a (non canonical) isomorphism

$$
H^{q}\left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}(U)\right) \simeq \operatorname{Hom}\left(H_{c}^{n-q}(U, \mathscr{A}), M\right) \oplus \operatorname{Ext}\left(H_{c}^{n-q+1}(U, \mathscr{A}), M\right)
$$

and in the special case $q=0$ a canonical isomorphism

$$
\begin{equation*}
H^{0}\left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}(U)\right)=\operatorname{Hom}\left(H_{c}^{n}(U, \mathscr{A}), M\right) . \tag{16.9}
\end{equation*}
$$

For an open subset $\Omega$ homeomorphic to $\mathbb{R}^{n}$, we have $\mathscr{A}_{\mid \Omega} \simeq L$. Proposition 14.16 and the exact sequence of the pair yield

$$
H_{c}^{q}(\Omega, L) \simeq H^{q}\left(S^{n},\{\infty\} ; L\right)=\left\{\begin{array}{lll}
L & \text { for } q=n \\
0 & \text { for } & q \neq n
\end{array}\right.
$$

If $\Omega \simeq \mathbb{R}^{n}$, we find

$$
H^{0}\left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}(\Omega)\right) \simeq \operatorname{Hom}(L, M), \quad H^{1}\left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}(\Omega)\right) \simeq \operatorname{Ext}(L, M)
$$

and $H^{q}(\mathscr{F} \mathscr{A}, M(\Omega))=0$ for $q \neq 0,1$. Consider open sets $V \subset \Omega$ where $V$ is a deformation retract of $\Omega$. Then the restriction map $H^{q}\left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}(\Omega)\right) \longrightarrow H^{q}\left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}(V)\right)$ is an isomorphism. Taking the direct limit over all such neighborhoods $V$ of a given point $x \in \Omega$, we see that $\mathscr{H}^{0}\left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}\right)$ and $\mathscr{H}^{1}\left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}\right)$ are locally constant sheaves of stalks Hom $(L, M)$ and $\operatorname{Ext}(L, M)$, and that $\mathscr{H}^{q}\left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}\right)=0$ for $q \neq 0,1$. When $\operatorname{Ext}(L, M)=0$, the complex $\mathscr{F}_{\mathscr{A}, M}^{\bullet}$ is thus a flabby resolution of $\mathscr{H}^{0}=\mathscr{H}^{0}\left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}\right)$ and we get isomorphisms

$$
\begin{align*}
& H^{q}\left(\mathscr{F}_{\mathscr{A}, M}(X)\right)=H^{q}\left(X, \mathscr{H}^{0}\right),  \tag{16.10}\\
& \mathscr{H}^{0}(U)=H^{0}\left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}(U)\right)=\operatorname{Hom}\left(H_{c}^{n}(U, \mathscr{A}), M\right) . \tag{16.11}
\end{align*}
$$

(16.12) Definition. The locally constant sheaf $\tau_{X}=\mathscr{H}^{0}(\mathscr{F} \mathbb{Z}, \mathbb{Z})$ of stalk $\mathbb{Z}$ defined by

$$
\tau_{X}(U)=\operatorname{Hom}_{\mathbb{Z}}\left(H_{c}^{n}(U, \mathbb{Z}), \mathbb{Z}\right)
$$

is called the orientation sheaf (or dualizing sheaf) of $X$.

This sheaf is given by a homomorphism $\pi_{1}(X) \longrightarrow\{1,-1\}$; it is not difficult to check that $\tau_{X}$ coincides with the trivial sheaf $\mathbb{Z}$ if and only if $X$ is orientable (cf. exercice 18.?). In general, $H_{c}^{n}(U, \mathscr{A})=H_{c}^{n}(U, \mathbb{Z}) \otimes_{\mathbb{Z}} \mathscr{A}(U)$ for any small open set $U$ on which $\mathscr{A}$ is trivial, thus

$$
\mathscr{H}^{0}\left(\mathscr{F}_{\mathscr{A}, M}^{\bullet}\right)=\tau_{X} \otimes_{\mathbb{Z}} \operatorname{Hom}(\mathscr{A}, M) .
$$

A combination of (16.7) and (16.10) then gives:
(16.13) Poincaré duality theorem. Let $X$ be a topological manifold, let $\mathcal{A}$ be a locally constant sheaf over $X$ of stalk $L$ and let $M$ be a $R$-module such that $\operatorname{Ext}(L, M)=0$. There is a split exact sequence

$$
\begin{aligned}
0 \longrightarrow \operatorname{Ext}\left(H_{c}^{n-q+1}(X, \mathscr{A}), M\right) \longrightarrow H^{q} & \left(X, \tau_{X} \otimes \operatorname{Hom}(\mathscr{A}, M)\right) \\
& \longrightarrow \operatorname{Hom}\left(H_{c}^{n-q}(X, \mathscr{A}), M\right) \longrightarrow 0 .
\end{aligned}
$$

In particular, if either $X$ is orientable or $R$ has characteristic 2, then

$$
\begin{array}{r}
0 \longrightarrow \operatorname{Ext}\left(H_{c}^{n-q+1}(X, R), R\right) \longrightarrow H^{q}(X, R) \longrightarrow \operatorname{Hom}\left(H_{c}^{n-q}(X, R), R\right) \\
\longrightarrow 0
\end{array}
$$

(16.14) Corollary. Let $X$ be a connected topological manifold, $n=\operatorname{dim} X$. Then for any $R$-module $L$
a) $H_{c}^{n}\left(X, \tau_{X} \otimes L\right) \simeq L$;
b) $H_{c}^{n}(X, L) \simeq L / 2 L$ if $X$ is not orientable.

Proof. First assume that $L$ is free. For $q=0$ and $\mathscr{A}=\tau_{X} \otimes L$, the Poincaré duality formula gives an isomorphism

$$
\operatorname{Hom}\left(H_{c}^{n}\left(X, \tau_{X} \otimes L\right), M\right) \simeq \operatorname{Hom}(L, M)
$$

and the isomorphism is functorial with respect to morphisms $M \longrightarrow M^{\prime}$. Taking $M=L$ or $M=H_{c}^{n}\left(X, \tau_{X} \otimes L\right)$, we easily obtain the existence of inverse morphisms

$$
H_{c}^{n}\left(X, \tau_{X} \otimes L\right) \longrightarrow L \quad \text { and } \quad L \longrightarrow H_{c}^{n}\left(X, \tau_{X} \otimes L\right)
$$

hence equality a). Similarly, for $\mathcal{A}=L$ we get

$$
\operatorname{Hom}\left(H_{c}^{n}(X, L), M\right) \simeq H^{0}\left(X, \tau_{X} \otimes \operatorname{Hom}(L, M)\right)
$$

If $X$ is non orientable, then $\tau_{X}$ is non trivial and the global sections of the sheaf $\tau_{X} \otimes$ $\operatorname{Hom}(L, M)$ consist of 2 -torsion elements of $\operatorname{Hom}(L, M)$, that is

$$
\operatorname{Hom}\left(H_{c}^{n}(X, L), M\right) \simeq \operatorname{Hom}(L / 2 L, M)
$$

Formula b) follows. If $L$ is not free, the result can be extended by using a free resolution $0 \rightarrow L_{1} \rightarrow L_{0} \rightarrow L \rightarrow 0$ and the associated long exact sequence.
(16.15) Remark. If $X$ is a connected non compact $n$-dimensional manifold, it can be proved (exercise 18.?) that $H^{n}(X, \mathscr{A})=0$ for every locally constant sheaf $\mathscr{A}$ on $X$.

Assume from now on that $X$ is oriented. Replacing $M$ by $L \otimes M$ and using the obvious morphism $M \longrightarrow \operatorname{Hom}(L, L \otimes M)$, the Poincaré duality theorem yields a morphism

$$
\begin{equation*}
H^{q}(X, M) \longrightarrow \operatorname{Hom}\left(H_{c}^{n-q}(X, L), L \otimes M\right), \tag{16.16}
\end{equation*}
$$

in other words, a bilinear pairing

$$
H_{c}^{n-q}(X, L) \times H^{q}(X, M) \longrightarrow L \otimes M .
$$

(16.17) Proposition. Up to the sign, the above pairing is given by the cup product, modulo the identification $H_{c}^{n}(X, L \otimes M) \simeq L \otimes M$.

Proof. By functoriality in $L$, we may assume $L=R$. Then we make the following special choices of resolutions:

$$
\begin{aligned}
\mathscr{L}^{q} & =R^{[q]} \quad \text { for } q<n, \quad \mathscr{L}^{n}=\operatorname{ker}\left(R^{[q]} \longrightarrow R^{[q+1]}\right), \\
M^{0} & =\text { an injective module containing } M_{c}^{[n]}(X) / d^{n-1} M_{c}^{[n-1]}(X) .
\end{aligned}
$$

We embed $M$ in $M^{0}$ by $\lambda \mapsto u \otimes_{\mathbb{Z}} \lambda$ where $u \in \mathbb{Z}^{[n]}(X)$ is a representative of a generator of $H_{c}^{n}(X, \mathbb{Z})$, and we set $M^{1}=M^{0} / M$. The projection map $M^{0} \longrightarrow M^{1}$ can be seen as an extension of

$$
\widetilde{d}^{n}: \quad M_{c}^{[n]}(X) / d^{n-1} M_{c}^{[n-1]}(X) \longrightarrow d^{n} M_{c}^{[n]}(X)
$$

since Ker $\widetilde{d}^{n} \simeq H_{c}^{n}(X, M)=M$. The inclusion $d^{n} M_{c}^{[n]}(X) \subset M^{1}$ can be extended into a map $\pi: M_{c}^{[n+1]}(X) \longrightarrow M^{1}$. The cup product bilinear map

$$
M^{[q]}(U) \times R_{c}^{[n-q]}(U) \longrightarrow M_{c}^{[n]}(X) \longrightarrow M^{0}
$$

gives rise to a morphism $M^{[q]}(U) \longrightarrow \mathscr{F}_{R, M}^{q}(U)$ defined by

$$
\begin{gather*}
M^{[q]}(U) \longrightarrow \operatorname{Hom}\left(\mathscr{L}_{c}^{n-q}(U), M^{0}\right) \oplus \operatorname{Hom}\left(\mathscr{L}_{c}^{n-q+1}(U), M^{1}\right) \\
f \longmapsto \quad(g \longmapsto f \smile g) \quad \oplus \quad(h \longmapsto \pi(f \smile h)) . \tag{16.18}
\end{gather*}
$$

This morphism is easily seen to give a morphism of differential sheaves $M^{[\bullet]} \longrightarrow \mathscr{F}_{R, M}^{\bullet}$, when $M^{[\bullet]}$ is truncated in degree $n$, i.e. when $M^{[n]}$ is replaced by $\operatorname{Ker} d^{n}$. The induced morphism

$$
M=\mathscr{H}^{0}\left(M^{[\bullet]}\right) \longrightarrow \mathscr{H}^{0}\left(\mathscr{F}_{R, M}^{\bullet}\right)
$$

is then the identity map, hence the cup product morphism (16.18) actually induces the Poincaré duality map (16.16).
(16.19) Remark. If $X$ is an oriented differentiable manifold, the natural isomorphism $H_{c}^{n}(X, \mathbb{R}) \simeq \mathbb{R}$ given by 16.14 a) corresponds in De Rham cohomology to the integration morphism $f \longmapsto \int_{X} f, f \in \mathscr{D}_{n}(X)$. Indeed, by a partition of unity, we may assume that Supp $f \subset \Omega \simeq \mathbb{R}^{n}$. The proof is thus reduced to the case $X=\mathbb{R}^{n}$, which itself reduces
to $X=\mathbb{R}$ since the cup product is compatible with the wedge product of forms. Let us consider the covering $\mathscr{U}=(] k-1, k+1[)_{k \in \mathbb{Z}}$ and a partition of unity $\left(\psi_{k}\right)$ subordinate to $\boldsymbol{U}$. The Čech differential

$$
\begin{aligned}
A C^{0}(\cup, \mathbb{U}) & \longrightarrow A C^{1}(\mathscr{U}, \mathbb{Z}) \\
\left(c_{k}\right) & \longmapsto\left(c_{k k+1}\right)=\left(c_{k+1}-c_{k}\right)
\end{aligned}
$$

shows immediately that the generators of $H_{c}^{1}(\mathbb{R}, \mathbb{Z})$ are the 1-cocycles $c$ such that $c_{01}=$ $\pm 1$ and $c_{k k+1}=0$ for $k \neq 0$. By Formula (6.12), the associated closed differential form is

$$
f=c_{01} \psi_{1} d \psi_{0}+c_{10} \psi_{0} d \psi_{1}
$$

hence $f= \pm \mathbf{1}_{[0,1]} d \psi_{0}$ and $f$ does satisfy $\int_{\mathbb{R}} f= \pm 1$.
(16.20) Corollary. If $X$ is an oriented $\mathscr{C}^{\infty}$ manifold, the bilinear map

$$
H_{c}^{n-q}(X, \mathbb{R}) \times H^{q}(X, \mathbb{R}) \longrightarrow \mathbb{R}, \quad(\{f\},\{g\}) \longmapsto \int_{X} f \wedge g
$$

is well defined and identifies $H^{q}(X, \mathbb{R})$ to the dual of $H_{c}^{n-q}(X, \mathbb{R})$.

## Chapter V

## Hermitian Vector Bundles

This chapter introduces the basic definitions concerning vector bundles and connections. In the first sections, the concepts of connection, curvature form, first Chern class are considered in the framework of differentiable manifolds. Although we are mainly interested in complex manifolds, the ideas which will be developed in the next chapters also involve real analysis and real geometry as essential tools. In the second part, the special features of connections over complex manifolds are investigated in detail: Chern connections, first Chern class of type ( 1,1 ), induced curvature forms on sub- and quotient bundles, .... These general concepts are then illustrated by the example of universal vector bundles over $\mathbb{P}^{n}$ and over Grassmannians.

## § 1. Definition of Vector Bundles

Let $M$ be a $\mathscr{C}^{\infty}$ differentiable manifold of dimension $m$ and let $\mathbb{K}=\mathbb{R}$ or $\mathbb{K}=\mathbb{C}$ be the scalar field. A (real, complex) vector bundle of rank $r$ over $M$ is a $\mathscr{C}^{\infty}$ manifold $E$ together with
i) a $\mathscr{C}^{\infty} \operatorname{map} \pi: E \longrightarrow M$ called the projection,
ii) a $\mathbb{K}$-vector space structure of dimension $r$ on each fiber $E_{x}=\pi^{-1}(x)$
such that the vector space structure is locally trivial. This means that there exists an open covering $\left(V_{\alpha}\right)_{\alpha \in I}$ of $M$ and $\mathscr{C}^{\infty}$ diffeomorphisms called trivializations

$$
\theta_{\alpha}: E_{\mid V_{\alpha}} \longrightarrow V_{\alpha} \times \mathbb{K}^{r}, \quad \text { where } \quad E_{\mid V_{\alpha}}=\pi^{-1}\left(V_{\alpha}\right),
$$

such that for every $x \in V_{\alpha}$ the map

$$
E_{x} \xrightarrow{\theta_{\alpha}}\{x\} \times \mathbb{K}^{r} \longrightarrow \mathbb{K}^{r}
$$

is a linear isomorphism. For each $\alpha, \beta \in I$, the map

$$
\theta_{\alpha \beta}=\theta_{\alpha} \circ \theta_{\beta}^{-1}:\left(V_{\alpha} \cap V_{\beta}\right) \times \mathbb{K}^{r} \longrightarrow\left(V_{\alpha} \cap V_{\beta}\right) \times \mathbb{K}^{r}
$$

acts as a linear automorphism on each fiber $\{x\} \times \mathbb{K}^{r}$. It can thus be written

$$
\theta_{\alpha \beta}(x, \xi)=\left(x, g_{\alpha \beta}(x) \cdot \xi\right), \quad(x, \xi) \in\left(V_{\alpha} \cap V_{\beta}\right) \times \mathbb{K}^{r}
$$

where $\left(g_{\alpha \beta}\right)_{(\alpha, \beta) \in I \times I}$ is a collection of invertible matrices with coefficients in $\mathscr{C}^{\infty}\left(V_{\alpha} \cap\right.$ $\left.V_{\beta}, \mathbb{K}\right)$, satisfying the cocycle relation

$$
\begin{equation*}
g_{\alpha \beta} g_{\beta \gamma}=g_{\alpha \gamma} \quad \text { on } \quad V_{\alpha} \cap V_{\beta} \cap V_{\gamma} . \tag{1.1}
\end{equation*}
$$

The collection $\left(g_{\alpha \beta}\right)$ is called a system of transition matrices. Conversely, any collection of invertible matrices satisfying (1.1) defines a vector bundle $E$, obtained by gluing the charts $V_{\alpha} \times \mathbb{K}^{r}$ via the identifications $\theta_{\alpha \beta}$.
(1.2) Example. The product manifold $E=M \times \mathbb{K}^{r}$ is a vector bundle over $M$, and is called the trivial vector bundle of rank $r$ over $M$. We shall often simply denote it $\mathbb{K}^{r}$ for brevity.
(1.3) Example. A much more interesting example of real vector bundle is the tangent bundle $T M$; if $\tau_{\alpha}: V_{\alpha} \longrightarrow \mathbb{R}^{n}$ is a collection of coordinate charts on $M$, then $\theta_{\alpha}=$ $\pi \times d \tau_{\alpha}: T M_{\mid V_{\alpha}} \longrightarrow V_{\alpha} \times \mathbb{R}^{m}$ define trivializations of $T M$ and the transition matrices are given by $g_{\alpha \beta}(x)=d \tau_{\alpha \beta}\left(x^{\beta}\right)$ where $\tau_{\alpha \beta}=\tau_{\alpha} \circ \tau_{\beta}^{-1}$ and $x^{\beta}=\tau_{\beta}(x)$. The dual $T^{\star} M$ of $T M$ is called the cotangent bundle and the $p$-th exterior power $\Lambda^{p} T^{\star} M$ is called the bundle of differential forms of degree $p$ on $M$.
(1.4) Definition. If $\Omega \subset M$ is an open subset and $k$ a positive integer or $+\infty$, we let $\mathscr{C}^{k}(\Omega, E)$ denote the space of $\mathscr{C}^{k}$ sections of $E_{\mid \Omega}$, i.e. the space of $\mathscr{C}^{k}$ maps $s: \Omega \longrightarrow E$ such that $s(x) \in E_{x}$ for all $x \in \Omega\left(\right.$ that is $\left.\pi \circ s=\operatorname{Id}_{\Omega}\right)$.

Let $\theta: E_{\mid V} \longrightarrow V \times \mathbb{K}^{r}$ be a trivialization of $E$. To $\theta$, we associate the $\mathscr{C}^{\infty}$ frame $\left(e_{1}, \ldots, e_{r}\right)$ of $E_{\mid V}$ defined by

$$
e_{\lambda}(x)=\theta^{-1}\left(x, \varepsilon_{\lambda}\right), \quad x \in V,
$$

where $\left(\varepsilon_{\lambda}\right)$ is the standard basis of $\mathbb{K}^{r}$. A section $s \in \mathscr{C}^{k}(V, E)$ can then be represented in terms of its components $\theta(s)=\sigma=\left(\sigma_{1}, \ldots, \sigma_{r}\right)$ by

$$
s=\sum_{1 \leqslant \lambda \leqslant r} \sigma_{\lambda} e_{\lambda} \quad \text { on } \quad V, \quad \sigma_{\lambda} \in \mathscr{C}^{k}(V, \mathbb{K})
$$

Let $\left(\theta_{\alpha}\right)$ be a family of trivializations relative to a covering $\left(V_{\alpha}\right)$ of $M$. Given a global section $s \in \mathscr{C}^{k}(M, E)$, the components $\theta_{\alpha}(s)=\sigma^{\alpha}=\left(\sigma_{1}^{\alpha}, \ldots, \sigma_{r}^{\alpha}\right)$ satisfy the transition relations

$$
\begin{equation*}
\sigma^{\alpha}=g_{\alpha \beta} \sigma^{\beta} \quad \text { on } \quad V_{\alpha} \cap V_{\beta} . \tag{1.5}
\end{equation*}
$$

Conversely, any collection of vector valued functions $\sigma^{\alpha}: V_{\alpha} \longrightarrow \mathbb{K}^{r}$ satisfying the transition relations defines a global section $s$ of $E$.

More generally, we shall also consider differential forms on $M$ with values in $E$. Such forms are nothing else than sections of the tensor product bundle $\Lambda^{p} T^{\star} M \otimes_{\mathbb{R}} E$. We shall write

$$
\begin{align*}
& \mathscr{C}_{p}^{k}(\Omega, E)=\mathscr{C}^{k}\left(\Omega, \Lambda^{p} T^{\star} M \otimes_{\mathbb{R}} E\right)  \tag{1.6}\\
& \mathscr{C}_{\bullet}^{k}(\Omega, E)=\bigoplus_{0 \leqslant p \leqslant m} \mathscr{C}_{p}^{k}(\Omega, E) \tag{1.7}
\end{align*}
$$

## § 2. Linear Connections

A (linear) connection $D$ on the bundle $E$ is a linear differential operator of order 1 acting on $\mathscr{C}_{\bullet}^{\infty}(M, E)$ and satisfying the following properties:

$$
\begin{align*}
& D: \mathscr{C}_{q}^{\infty}(M, E) \longrightarrow \mathscr{C}_{q+1}^{\infty}(M, E),  \tag{2.1}\\
& D(f \wedge s)=d f \wedge s+(-1)^{p} f \wedge D s
\end{align*}
$$

for any $f \in \mathscr{C}_{p}^{\infty}(M, \mathbb{K})$ and $s \in \mathscr{C}_{q}^{\infty}(M, E)$, where $d f$ stands for the usual exterior derivative of $f$.

Assume that $\theta: E_{\upharpoonright \Omega} \rightarrow \Omega \times \mathbb{K}^{r}$ is a trivialization of $E_{\upharpoonright \Omega}$, and let $\left(e_{1}, \ldots, e_{r}\right)$ be the corresponding frame of $E_{\uparrow \Omega}$. Then any $s \in \mathscr{C}_{q}^{\infty}(\Omega, E)$ can be written in a unique way

$$
s=\sum_{1 \leqslant \lambda \leqslant r} \sigma_{\lambda} \otimes e_{\lambda}, \quad \sigma_{\lambda} \in \mathscr{C}_{q}^{\infty}(\Omega, \mathbb{K}) .
$$

By axiom (2.1') we get

$$
D s=\sum_{1 \leqslant \lambda \leqslant r}\left(d \sigma_{\lambda} \otimes e_{\lambda}+(-1)^{p} \sigma_{\lambda} \wedge D e_{\lambda}\right) .
$$

If we write $D e_{\mu}=\sum_{1 \leqslant \lambda \leqslant r} a_{\lambda \mu} \otimes e_{\lambda}$ where $a_{\lambda \mu} \in \mathscr{C}_{1}^{\infty}(\Omega, \mathbb{K})$, we thus have

$$
D s=\sum_{\lambda}\left(d \sigma_{\lambda}+\sum_{\mu} a_{\lambda \mu} \wedge \sigma_{\mu}\right) \otimes e_{\lambda} .
$$

Identify $E_{\lceil\Omega}$ with $\Omega \times \mathbb{K}^{r}$ via $\theta$ and denote by $d$ the trivial connection $d \sigma=\left(d \sigma_{\lambda}\right)$ on $\Omega \times \mathbb{K}^{r}$. Then the operator $D$ can be written

$$
\begin{equation*}
D s \simeq_{\theta} d \sigma+A \wedge \sigma \tag{2.2}
\end{equation*}
$$

where $A=\left(a_{\lambda \mu}\right) \in \mathscr{C}_{1}^{\infty}\left(\Omega, \operatorname{Hom}\left(\mathbb{K}^{r}, \mathbb{K}^{r}\right)\right)$. Conversely, it is clear that any operator $D$ defined in such a way is a connection on $E_{\mid \Omega}$. The matrix 1-form $A$ will be called the connection form of $D$ associated to the trivialization $\theta$. If $\widetilde{\theta}: E_{\lceil\Omega} \rightarrow \Omega \times \mathbb{K}^{r}$ is another trivialization and if we set

$$
g=\widetilde{\theta} \circ \theta^{-1} \in \mathscr{C}^{\infty}\left(\Omega, \operatorname{Gl}\left(\mathbb{K}^{r}\right)\right)
$$

then the new components $\widetilde{\sigma}=\left(\widetilde{\sigma}_{\lambda}\right)$ are related to the old ones by $\widetilde{\sigma}=g \sigma$. Let $\widetilde{A}$ be the connection form of $D$ with respect to $\widetilde{\theta}$. Then

$$
\begin{aligned}
D s & \simeq_{\widetilde{\theta}} d \widetilde{\sigma}+\widetilde{A} \wedge \widetilde{\sigma} \\
D s & \simeq_{\theta} g^{-1}(d \widetilde{\sigma}+\widetilde{A} \wedge \widetilde{\sigma})=g^{-1}(d(g \sigma)+\widetilde{A} \wedge g \sigma) \\
& =d \sigma+\left(g^{-1} \widetilde{A} g+g^{-1} d g\right) \wedge \sigma .
\end{aligned}
$$

Therefore we obtain the gauge transformation law:

$$
\begin{equation*}
A=g^{-1} \widetilde{A} g+g^{-1} d g \tag{2.3}
\end{equation*}
$$

## §3. Curvature Tensor

Let us compute $D^{2}: \mathscr{C}_{q}^{\infty}(M, E) \rightarrow \mathscr{C}_{q+2}^{\infty}(M, E)$ with respect to the trivialization $\theta: E_{\lceil\Omega} \rightarrow \Omega \times \mathbb{K}^{r}$. We obtain

$$
\begin{aligned}
D^{2} s & \simeq_{\theta} d(d \sigma+A \wedge \sigma)+A \wedge(d \sigma+A \wedge \sigma) \\
& =d^{2} \sigma+(d A \wedge \sigma-A \wedge d \sigma)+(A \wedge d \sigma+A \wedge A \wedge \sigma) \\
& =(d A+A \wedge A) \wedge \sigma
\end{aligned}
$$

It follows that there exists a global 2-form $\Theta(D) \in \mathscr{C}_{2}^{\infty}(M, \operatorname{Hom}(E, E))$ called the curvature tensor of $D$, such that

$$
D^{2} s=\Theta(D) \wedge s
$$

given with respect to any trivialization $\theta$ by

$$
\begin{equation*}
\Theta(D) \simeq_{\theta} d A+A \wedge A \tag{3.1}
\end{equation*}
$$

(3.2) Remark. If $E$ is of $\operatorname{rank} r=1$, then $A \in \mathscr{C}_{1}^{\infty}(M, \mathbb{K})$ and $\operatorname{Hom}(E, E)$ is canonically isomorphic to the trivial bundle $M \times \mathbb{K}$, because the endomorphisms of each fiber $E_{x}$ are homotheties. With the identification $\operatorname{Hom}(E, E)=\mathbb{K}$, the curvature tensor $\Theta(D)$ can be considered as a closed 2 -form with values in $\mathbb{K}$ :

$$
\begin{equation*}
\Theta(D)=d A \tag{3.3}
\end{equation*}
$$

In this case, the gauge transformation law can be written

$$
\begin{equation*}
A=\widetilde{A}+g^{-1} d g, \quad g=\widetilde{\theta} \circ \theta^{-1} \in \mathscr{C}^{\infty}\left(\Omega, \mathbb{K}^{\star}\right) \tag{3.4}
\end{equation*}
$$

It is then immediately clear that $d A=d \widetilde{A}$, and this equality shows again that $d A$ does not depend on $\theta$.

Now, we show that the curvature tensor is closely related to commutation properties of covariant derivatives.
(3.5) Definition. If $\xi$ is a $\mathscr{C}^{\infty}$ vector field with values in $T M$, the covariant derivative of a section $s \in \mathscr{C}^{\infty}(M, E)$ in the direction $\xi$ is the section $\xi_{D} \cdot s \in \mathscr{C}^{\infty}(M, E)$ defined by $\xi_{D} \cdot s=D s \cdot \xi$.
(3.6) Proposition. For all sections $s \in \mathscr{C}^{\infty}(M, E)$ and all vector fields $\xi, \eta \in$ $\complement^{\infty}(M, T M)$, we have

$$
\xi_{D} \cdot\left(\eta_{D} \cdot s\right)-\eta_{D} \cdot\left(\xi_{D} \cdot s\right)=[\xi, \eta]_{D} \cdot s+\Theta(D)(\xi, \eta) \cdot s
$$

where $[\xi, \eta] \in \mathscr{C}^{\infty}(M, T M)$ is the Lie bracket of $\xi, \eta$.
Proof. Let $\left(x_{1}, \ldots, x_{m}\right)$ be local coordinates on an open set $\Omega \subset M$. Let $\theta: E_{\upharpoonright \Omega} \longrightarrow$ $\Omega \times \mathbb{K}^{r}$ be a trivialization of $E$ and let $A$ be the corresponding connection form. If $\xi=\sum \xi_{j} \partial / \partial x_{j}$ and $A=\sum A_{j} d x_{j}$, we find

$$
\begin{equation*}
\xi_{D} s \simeq_{\theta}(d \sigma+A \sigma) \cdot \xi=\sum_{j} \xi_{j}\left(\frac{\partial \sigma}{\partial x_{j}}+A_{j} \cdot \sigma\right) . \tag{3.7}
\end{equation*}
$$

Now, we compute the above commutator $\left[\xi_{D}, \eta_{D}\right]$ at a given point $z_{0} \in \Omega$. Without loss of generality, we may assume $A\left(z_{0}\right)=0$; in fact, one can always find a gauge transformation $g$ near $z_{0}$ such that $g\left(z_{0}\right)=\operatorname{Id}$ and $d g\left(z_{0}\right)=A\left(z_{0}\right)$; then (2.3) yields $\widetilde{A}\left(z_{0}\right)=0$. If $\eta=\sum \eta_{k} \partial / \partial x_{k}$, we find $\eta_{D} \cdot s \simeq_{\theta} \sum \eta_{k} \partial \sigma / \partial x_{k}$ at $z_{0}$, hence

$$
\begin{aligned}
\eta_{D} \cdot\left(\xi_{D} \cdot s\right) & \simeq_{\theta} \sum_{k} \eta_{k} \frac{\partial}{\partial x_{k}} \sum_{j} \xi_{j}\left(\frac{\partial \sigma}{\partial x_{j}}+A_{j} \cdot \sigma\right), \\
\xi_{D} \cdot\left(\eta_{D} \cdot s\right) & -\eta_{D} \cdot\left(\xi_{D} \cdot s\right) \simeq_{\theta} \\
& \simeq_{\theta} \sum_{j, k}\left(\xi_{k} \frac{\partial \eta_{j}}{\partial x_{k}}-\eta_{k} \frac{\partial \xi_{j}}{\partial x_{k}}\right) \frac{\partial \sigma}{\partial x_{j}}+\sum_{j, k} \frac{\partial A_{j}}{\partial x_{k}}\left(\xi_{j} \eta_{k}-\eta_{j} \xi_{k}\right) \cdot \sigma \\
& =d \sigma([\xi, \eta])+d A(\xi, \eta) \cdot \sigma,
\end{aligned}
$$

whereas $\Theta(D) \simeq_{\theta} d A$ and $[\xi, \eta]_{D} s \simeq_{\theta} d \sigma([\xi, \eta])$ at point $z_{0}$.

## $\S$ 4. Operations on Vector Bundles

Let $E, F$ be vector bundles of rank $r_{1}, r_{2}$ over $M$. Given any functorial operation on vector spaces, a corresponding operation can be defined on bundles by applying the operation on each fiber. For example $E^{\star}, E \oplus F, \operatorname{Hom}(E, F)$ are defined by

$$
\left(E^{\star}\right)_{x}=\left(E_{x}\right)^{\star}, \quad(E \oplus F)_{x}=E_{x} \oplus F_{x}, \quad \operatorname{Hom}(E, F)_{x}=\operatorname{Hom}\left(E_{x}, F_{x}\right)
$$

The bundles $E$ and $F$ can be trivialized over the same covering $V_{\alpha}$ of $M$ (otherwise take a common refinement). If $\left(g_{\alpha \beta}\right)$ and $\left(\gamma_{\alpha \beta}\right)$ are the transition matrices of $E$ and $F$, then for example $E \otimes F, \Lambda^{k} E, E^{\star}$ are the bundles defined by the transition matrices $g_{\alpha \beta} \otimes \gamma_{\alpha \beta}$, $\Lambda^{k} g_{\alpha \beta},\left(g_{\alpha \beta}^{\dagger}\right)^{-1}$ where $\dagger$ denotes transposition.

Suppose now that $E, F$ are equipped with connections $D_{E}, D_{F}$. Then natural connections can be associated to all derived bundles. Let us mention a few cases. First, we let

$$
\begin{equation*}
D_{E \oplus F}=D_{E} \oplus D_{F} \tag{4.1}
\end{equation*}
$$

It follows immediately that

$$
\Theta\left(D_{E \oplus F}\right)=\Theta\left(D_{E}\right) \oplus \Theta\left(D_{F}\right)
$$

$D_{E \otimes F}$ will be defined in such a way that the usual formula for the differentiation of a product remains valid. For every $s \in \mathscr{C}_{\bullet}^{\infty}(M, E), t \in \mathscr{C}_{\bullet}^{\infty}(M, F)$, the wedge product $s \wedge t$ can be combined with the bilinear map $E \times F \longrightarrow E \otimes F$ in order to obtain a section $s \wedge t \in \mathscr{C}^{\infty}(M, E \otimes F)$ of degree $\operatorname{deg} s+\operatorname{deg} t$. Then there exists a unique connection $D_{E \otimes F}$ such that

$$
\begin{equation*}
D_{E \otimes F}(s \wedge t)=D_{E} s \wedge t+(-1)^{\operatorname{deg} s} s \wedge D_{F} t \tag{4.2}
\end{equation*}
$$

As the products $s \wedge t$ generate $\mathscr{C}_{\bullet}^{\infty}(M, E \otimes F)$, the uniqueness is clear. If $E, F$ are trivial on an open set $\Omega \subset M$ and if $A_{E}, A_{F}$, are their connection 1-forms, the induced
connection $D_{E \otimes F}$ is given by the connection form $A_{E} \otimes \operatorname{Id}_{F}+\operatorname{Id}_{E} \otimes A_{F}$. The existence follows. An easy computation shows that $D_{E \otimes F}^{2}(s \wedge t)=D_{E}^{2} s \wedge t+s \wedge D_{F}^{2} t$, thus

$$
\Theta\left(D_{E \otimes F}\right)=\Theta\left(D_{E}\right) \otimes \operatorname{Id}_{F}+\operatorname{Id}_{E} \otimes \Theta\left(D_{F}\right)
$$

Similarly, there are unique connections $D_{E^{\star}}$ and $D_{\operatorname{Hom}(E, F)}$ such that

$$
\begin{align*}
\left(D_{E^{\star}} u\right) \cdot s & =d(u \cdot s)-(-1)^{\operatorname{deg} u^{u} u \cdot D_{E} s},  \tag{4.3}\\
\left(D_{\mathrm{Hom}(E, F)} v\right) \cdot s & =D_{F}(v \cdot s)-(-1)^{\operatorname{deg} v} v \cdot D_{E} s \tag{4.4}
\end{align*}
$$

whenever $s \in \mathscr{C}_{\bullet}^{\infty}(M, E), u \in \mathscr{C}_{\bullet}^{\infty}\left(M, E^{\star}\right), v \in \mathscr{C}_{\bullet}^{\infty}(\operatorname{Hom}(E, F))$. It follows that

$$
0=d^{2}(u \cdot s)=\left(\Theta\left(D_{E^{\star}}\right) \cdot u\right) \cdot s+u \cdot\left(\Theta\left(D_{E}\right) \cdot s\right)
$$

If $\dagger$ denotes the transposition operator $\operatorname{Hom}(E, E) \rightarrow \operatorname{Hom}\left(E^{\star}, E^{\star}\right)$, we thus get

$$
\Theta\left(D_{E^{\star}}\right)=-\Theta\left(D_{E}\right)^{\dagger} .
$$

With the identification $\operatorname{Hom}(E, F)=E^{\star} \otimes F$, Formula (4.2') implies

$$
\Theta\left(D_{\operatorname{Hom}(E, F)}\right)=\operatorname{Id}_{E^{\star}} \otimes \Theta\left(D_{F}\right)-\Theta\left(D_{E}\right)^{\dagger} \otimes \operatorname{Id}_{F}
$$

Finally, $\Lambda^{k} E$ carries a natural connection $D_{\Lambda^{k} E}$. For every $s_{1}, \ldots, s_{k}$ in $\mathscr{C}_{\bullet}^{\infty}(M, E)$ of respective degrees $p_{1}, \ldots, p_{k}$, this connection satisfies

$$
\begin{align*}
& D_{\Lambda^{k} E}\left(s_{1} \wedge \ldots \wedge s_{k}\right)=\sum_{1 \leqslant j \leqslant k}(-1)^{p_{1}+\ldots+p_{j-1}} s_{1} \wedge \ldots D_{E} s_{j} \ldots \wedge s_{k},  \tag{4.5}\\
& \Theta\left(D_{\Lambda^{k} E}\right) \cdot\left(s_{1} \wedge \ldots \wedge s_{k}\right)=\sum_{1 \leqslant j \leqslant k} s_{1} \wedge \ldots \wedge \Theta\left(D_{E}\right) \cdot s_{j} \wedge \ldots \wedge s_{k} .
\end{align*}
$$

In particular, the determinant bundle, defined by $\operatorname{det} E=\Lambda^{r} E$ where $r$ is the rank of $E$, has a curvature form given by

$$
\begin{equation*}
\Theta\left(D_{\operatorname{det} E}\right)=\mathscr{T}_{E}\left(\Theta\left(D_{E}\right)\right) \tag{4.6}
\end{equation*}
$$

where $\mathscr{T}_{E}: \operatorname{Hom}(E, E) \longrightarrow \mathbb{K}$ is the trace operator. As a conclusion of this paragraph, we mention the following simple identity.
(4.7) Bianchi identity. $D_{\text {Hom }(E, E)}\left(\Theta\left(D_{E}\right)\right)=0$.

Proof. By definition of $D_{\operatorname{Hom}(E, E)}$, we find for any $s \in \mathscr{C}^{\infty}(M, E)$

$$
\begin{aligned}
D_{\operatorname{Hom}(E, E)}\left(\Theta\left(D_{E}\right)\right) \cdot s & =D_{E}\left(\Theta\left(D_{E}\right) \cdot s\right)-\Theta\left(D_{E}\right) \cdot\left(D_{E} s\right) \\
& =D_{E}^{3} s-D_{E}^{3} s=0 .
\end{aligned}
$$

## § 5. Pull-Back of a Vector Bundle

Let $\widetilde{M}, M$ be $\mathscr{C}^{\infty}$ manifolds and $\psi: \widetilde{M} \rightarrow M$ a smooth map. If $E$ is a vector bundle on $M$, one can define in a natural way a $\mathscr{C}^{\infty}$ vector bundle $\widetilde{\pi}: \widetilde{E} \rightarrow \widetilde{M}$ and a $\mathscr{C}^{\infty}$ linear morphism $\Psi: \widetilde{E} \rightarrow E$ such that the diagram

$$
\begin{array}{lll}
\widetilde{E} & \xrightarrow{\Psi} & E \\
\downarrow \tilde{\pi} & & \downarrow \pi \\
\widetilde{M} & \xrightarrow{\psi} & M
\end{array}
$$

commutes and such that $\Psi: \widetilde{E}_{x} \longrightarrow E_{\psi(x)}$ is an isomorphism for every $x \in M$. The bundle $\widetilde{E}$ can be defined by

$$
\begin{equation*}
\widetilde{E}=\{(\widetilde{x}, \xi) \in \widetilde{M} \times E ; \psi(\widetilde{x})=\pi(\xi)\} \tag{5.1}
\end{equation*}
$$

and the maps $\widetilde{\pi}$ and $\Psi$ are then the restrictions to $\widetilde{E}$ of the projections of $\widetilde{M} \times E$ on $\widetilde{M}$ and $E$ respectively.

If $\theta_{\alpha}: E_{\mid V_{\alpha}} \longrightarrow V_{\alpha} \times \mathbb{K}^{r}$ are trivializations of $E$, the maps

$$
\widetilde{\theta}_{\alpha}=\theta_{\alpha} \circ \Psi: \widetilde{E}_{\upharpoonright \psi^{-1}\left(V_{\alpha}\right)} \longrightarrow \psi^{-1}\left(V_{\alpha}\right) \times \mathbb{K}^{r}
$$

define trivializations of $\widetilde{E}$ with respect to the covering $\widetilde{V}_{\alpha}=\psi^{-1}\left(V_{\alpha}\right)$ of $\widetilde{M}$. The corresponding system of transition matrices is given by

$$
\begin{equation*}
\widetilde{g}_{\alpha \beta}=g_{\alpha \beta} \circ \psi \quad \text { on } \quad \widetilde{V}_{\alpha} \cap \widetilde{V}_{\beta} . \tag{5.2}
\end{equation*}
$$

(5.3) Definition. $\widetilde{E}$ is termed the pull-back of $E$ under the map $\psi$ and is denoted $\widetilde{E}=\psi^{\star} E$.

Let $D$ be a connection on $E$. If $\left(A_{\alpha}\right)$ is the collection of connection forms of $D$ with respect to the $\theta_{\alpha}$ 's, one can define a connection $\widetilde{D}$ on $\widetilde{E}$ by the collection of connection forms $\widetilde{A}_{\alpha}=\psi^{\star} A_{\alpha} \in \mathscr{C}_{1}^{\infty}\left(\widetilde{V}_{\alpha}, \operatorname{Hom}\left(\mathbb{K}^{r}, \mathbb{K}^{r}\right)\right)$, i.e. for every $\widetilde{s} \in \mathscr{C}_{p}^{\infty}\left(\widetilde{V}_{\alpha}, \widetilde{E}\right)$

$$
\widetilde{D} \widetilde{s} \simeq_{\widetilde{\theta}_{\alpha}} d \widetilde{\sigma}+\psi^{\star} A_{\alpha} \wedge \widetilde{\sigma}
$$

Given any section $s \in \mathscr{C}_{p}^{\infty}(M, E)$, one defines a pull back $\psi^{\star} s$ which is a section in $\mathscr{C}_{p}^{\infty}(\widetilde{M}, \widetilde{E}):$ for $s=f \otimes u, f \in \mathscr{C}_{p}^{\infty}(M, \mathbb{K}), u \in \mathscr{C}^{\infty}(M, E)$, set $\psi^{\star} s=\psi^{\star} f \otimes(u \circ \psi)$. Then we have the formula

$$
\begin{equation*}
\widetilde{D}\left(\psi^{\star} s\right)=\psi^{\star}(D s) . \tag{5.4}
\end{equation*}
$$

Using (5.4), a simple computation yields

$$
\begin{equation*}
\Theta(\widetilde{D})=\psi^{\star}(\Theta(D)) \tag{5.5}
\end{equation*}
$$

## § 6. Parallel Translation and Flat Vector Bundles

Let $\gamma:[0,1] \longrightarrow M$ be a smooth curve and $s:[0,1] \rightarrow E$ a $\mathscr{C}^{\infty}$ section of $E$ along $\gamma$, i.e. a $\mathscr{C}_{\sim}^{\infty}$ map $s$ such that $s(t) \in E_{\gamma(t)}$ for all $t \in[0,1]$. Then $s$ can be viewed as a section of $\widetilde{E}=\gamma^{\star} E$ over $[0,1]$. The covariant derivative of $s$ is the section of $E$ along $\gamma$ defined by

$$
\begin{equation*}
\frac{D s}{d t}=\widetilde{D} s(t) \cdot \frac{d}{d t} \in E_{\gamma(t)} \tag{6.1}
\end{equation*}
$$

where $\widetilde{D}$ is the induced connection on $\widetilde{E}$. If $A$ is a connection form of $D$ with respect to a trivialization $\theta: E_{\mid \Omega} \longrightarrow \Omega \times \mathbb{K}^{r}$, we have $\widetilde{D} s \simeq_{\theta} d \sigma+\gamma^{\star} A \cdot \sigma$, i.e.

$$
\begin{equation*}
\frac{D s}{d t} \simeq_{\theta} \frac{d \sigma}{d t}+\left(A(\gamma(t)) \cdot \gamma^{\prime}(t)\right) \cdot \sigma(t) \quad \text { for } \quad \gamma(t) \in \Omega \tag{6.2}
\end{equation*}
$$

For $v \in E_{\gamma(0)}$ given, the Cauchy uniqueness and existence theorem for ordinary linear differential equations implies that there exists a unique section $s$ of $\widetilde{E}$ such that $s(0)=v$ and $D s / d t=0$.
(6.3) Definition. The linear map

$$
T_{\gamma}: E_{\gamma(0)} \longrightarrow E_{\gamma(1)}, \quad v=s(0) \longmapsto s(1)
$$

is called parallel translation along $\gamma$.

If $\gamma=\gamma_{2} \gamma_{1}$ is the composite of two paths $\gamma_{1}, \gamma_{2}$ such that $\gamma_{2}(0)=\gamma_{1}(1)$, it is clear that $T_{\gamma}=T_{\gamma_{2}} \circ T_{\gamma_{1}}$, and the inverse path $\gamma^{-1}: t \mapsto \gamma(1-t)$ is such that $T_{\gamma^{-1}}=T_{\gamma}^{-1}$. It follows that $T_{\gamma}$ is a linear isomorphism from $E_{\gamma(0)}$ onto $E_{\gamma(1)}$.

More generally, if $h: W \longrightarrow M$ is a $\mathscr{C}^{\infty}$ map from a domain $W \subset \mathbb{R}^{p}$ into $M$ and if $s$ is a section of $h^{\star} E$, we define covariant derivatives $D s / \partial t_{j}, 1 \leqslant j \leqslant p$, by $\widetilde{D}=h^{\star} D$ and

$$
\begin{equation*}
\frac{D s}{\partial t_{j}}=\widetilde{D} s \cdot \frac{\partial}{\partial t_{j}} . \tag{6.4}
\end{equation*}
$$

Since $\partial / \partial t_{j}, \partial / \partial t_{k}$ commute and since $\Theta(\widetilde{D})=h^{\star} \Theta(D)$, Prop. 3.6 implies

$$
\begin{equation*}
\frac{D}{\partial t_{j}} \frac{D s}{\partial t_{k}}-\frac{D}{\partial t_{k}} \frac{D s}{\partial t_{j}}=\Theta(\widetilde{D})\left(\frac{\partial}{\partial t_{j}}, \frac{\partial}{\partial t_{k}}\right) \cdot s=\Theta(D)_{h(t)}\left(\frac{\partial h}{\partial t_{j}}, \frac{\partial h}{\partial t_{k}}\right) \cdot s(t) \tag{6.5}
\end{equation*}
$$

(6.6) Definition. The connection $D$ is said to be flat if $\Theta(D)=0$.

Assume from now on that $D$ is flat. We then show that $T_{\gamma}$ only depends on the homotopy class of $\gamma$. Let $h:[0,1] \times[0,1] \longrightarrow M$ be a smooth homotopy $h(t, u)=\gamma_{u}(t)$ from $\gamma_{0}$ to $\gamma_{1}$ with fixed end points $a=\gamma_{u}(0), b=\gamma_{u}(1)$. Let $v \in E_{a}$ be given and let $s(t, u)$ be such that $s(0, u)=v$ and $D s / \partial t=0$ for all $u \in[0,1]$. Then $s$ is $\mathscr{C}^{\infty}$ in both variables $(t, u)$ by standard theorems on the dependence of parameters. Moreover (6.5) implies that the covariant derivatives $D / \partial t, D / \partial u$ commute. Therefore, if we set $s^{\prime}=D s / \partial u$, we find $D s^{\prime} / \partial t=0$ with initial condition $s^{\prime}(0, u)=0$ (recall that $s(0, u)$
is a constant). The uniqueness of solutions of differential equations implies that $s^{\prime}$ is identically zero on $[0,1] \times[0,1]$, in particular $T_{\gamma_{u}}(v)=s(1, u)$ must be constant, as desired.
(6.7) Proposition. Assume that $D$ is flat. If $\Omega$ is a simply connected open subset of $M$, then $E_{\uparrow \Omega}$ admits a $\mathscr{C}^{\infty}$ parallel frame $\left(e_{1}, \ldots, e_{r}\right)$, in the sense that $D e_{\lambda}=0$ on $\Omega$, $1 \leqslant \lambda \leqslant r$. For any two simply connected open subsets $\Omega, \Omega^{\prime}$ the transition automorphism between the corresponding parallel frames $\left(e_{\lambda}\right)$ and $\left(e_{\lambda}^{\prime}\right)$ is locally constant.

The converse statement " $E$ has parallel frames near every point implies that $\Theta(D)=$ 0 " can be immediately verified from the equality $\Theta(D)=D^{2}$.

Proof. Choose a base point $a \in \Omega$ and define a linear isomorphism $\Phi: \Omega \times E_{a} \longrightarrow E_{\uparrow \Omega}$ by sending $(x, v)$ on $T_{\gamma}(v) \in E_{x}$, where $\gamma$ is any path from $a$ to $x$ in $\Omega$ (two such paths are always homotopic by hypothesis). Now, for any path $\gamma$ from $a$ to $x$, we have by construction $(D / d t) \Phi(\gamma(t), v)=0$. Set $e_{v}(x)=\Phi(x, v)$. As $\gamma$ may reach any point $x \in \Omega$ with an arbitrary tangent vector $\xi=\gamma^{\prime}(1) \in T_{x} M$, we get $D e_{v}(x) \cdot \xi=$ $(D / d t) \Phi(\gamma(t), v)_{\mid t=1}=0$. Hence $D e_{v}$ is parallel for any fixed vector $v \in E_{a} ;$ Prop. 6.7 follows.

Assume that $M$ is connected. Let $a$ be a base point and $\widetilde{M} \longrightarrow M$ the universal covering of $M$. The manifold $\widetilde{M}$ can be considered as the set of pairs $(x,[\gamma])$, where $[\gamma]$ is a homotopy class of paths from $a$ to $x$. Let $\pi_{1}(M)$ be the fundamental group of $M$ with base point $a$, acting on $\widetilde{M}$ on the left by $[\kappa] \cdot(x,[\gamma])=\left(x,\left[\gamma \kappa^{-1}\right]\right)$. If $D$ is flat, $\pi_{1}(M)$ acts also on $E_{a}$ by $([\kappa], v) \mapsto T_{\kappa}(v),[\kappa] \in \pi_{1}(M), v \in E_{a}$, and we have a well defined map

$$
\Psi: \widetilde{M} \times E_{a} \longrightarrow E, \quad \Psi(x,[\gamma])=T_{\gamma}(v)
$$

Then $\Psi$ is invariant under the left action of $\pi_{1}(M)$ on $\widetilde{M} \times E_{a}$ defined by

$$
[\kappa] \cdot((x,[\gamma]), v)=\left(\left(x,\left[\gamma \kappa^{-1}\right]\right), T_{\kappa}(v)\right),
$$

therefore we have an isomorphism $E \simeq\left(\widetilde{M} \times E_{a}\right) / \pi_{1}(M)$.
Conversely, let $S$ be a $\mathbb{K}$-vector space of dimension $r$ together with a left action of $\pi_{1}(M)$. The quotient $E=(\widetilde{M} \times S) / \pi_{1}(M)$ is a vector bundle over $M$ with locally constant transition automorphisms $\left(g_{\alpha \beta}\right)$ relatively to any covering $\left(V_{\alpha}\right)$ of $M$ by simply connected open sets. The relation $\sigma^{\alpha}=g_{\alpha \beta} \sigma^{\beta}$ implies $d \sigma^{\alpha}=g_{\alpha \beta} d \sigma^{\beta}$ on $V_{\alpha} \cap V_{\beta}$. We may therefore define a connection $D$ on $E$ by letting $D s \simeq_{\theta_{\alpha}} d \sigma^{\alpha}$ on each $V_{\alpha}$. Then clearly $\Theta(D)=0$.

## $\S$ 7. Hermitian Vector Bundles and Connections

A complex vector bundle $E$ is said to be hermitian if a positive definite hermitian form $\left.\left|\left.\right|^{2}\right.$ is given on each fiber $E_{x}$ in such a way that the map $\left.E \rightarrow \mathbb{R}_{+}, \xi \mapsto\right| \xi\right|^{2}$ is smooth. The notion of a euclidean (real) vector bundle is similar, so we leave the reader adapt our notations to that case.

Let $\theta: E_{\mid \Omega} \longrightarrow \Omega \times \mathbb{C}^{r}$ be a trivialization and let $\left(e_{1}, \ldots, e_{r}\right)$ be the corresponding frame of $E_{\lceil\Omega}$. The associated inner product of $E$ is given by a positive definite hermitian
matrix ( $h_{\lambda \mu}$ ) with $\mathscr{C}^{\infty}$ coefficients on $\Omega$, such that

$$
\left\langle e_{\lambda}(x), e_{\mu}(x)\right\rangle=h_{\lambda \mu}(x), \quad \forall x \in \Omega
$$

When $E$ is hermitian, one can define a natural sesquilinear map

$$
\begin{align*}
\mathscr{C}_{p}^{\infty}(M, E) \times \mathscr{C}_{q}^{\infty}(M, E) & \longrightarrow \mathscr{C}_{p+q}^{\infty}(M, \mathbb{C}) \\
(s, t) & \longmapsto\{s, t\} \tag{7.1}
\end{align*}
$$

combining the wedge product of forms with the hermitian metric on $E$; if $s=\sum \sigma_{\lambda} \otimes e_{\lambda}$, $t=\sum \tau_{\mu} \otimes e_{\mu}$, we let

$$
\{s, t\}=\sum_{1 \leqslant \lambda, \mu \leqslant r} \sigma_{\lambda} \wedge \bar{\tau}_{\mu}\left\langle e_{\lambda}, e_{\mu}\right\rangle
$$

A connection $D$ is said to be compatible with the hermitian structure of $E$, or briefly hermitian, if for every $s \in \mathscr{C}_{p}^{\infty}(M, E), t \in \mathscr{C}_{q}^{\infty}(M, E)$ we have

$$
\begin{equation*}
d\{s, t\}=\{D s, t\}+(-1)^{p}\{s, D t\} . \tag{7.2}
\end{equation*}
$$

Let $\left(e_{1}, \ldots, e_{r}\right)$ be an orthonormal frame of $E_{\lceil\Omega}$. Denote $\theta(s)=\sigma=\left(\sigma_{\lambda}\right)$ and $\theta(t)=$ $\tau=\left(\tau_{\lambda}\right)$. Then

$$
\begin{aligned}
\{s, t\} & =\{\sigma, \tau\}=\sum_{1 \leqslant \lambda \leqslant r} \sigma_{\lambda} \wedge \bar{\tau}_{\lambda}, \\
d\{s, t\} & =\{d \sigma, \tau\}+(-1)^{p}\{\sigma, d \tau\} .
\end{aligned}
$$

Therefore $D_{\lceil\Omega}$ is hermitian if and only if its connection form $A$ satisfies

$$
\{A \sigma, \tau\}+(-1)^{p}\{\sigma, A \tau\}=\left\{\left(A+A^{\star}\right) \wedge \sigma, \tau\right\}=0
$$

for all $\sigma, \tau$, i.e.

$$
\begin{equation*}
A^{\star}=-A \quad \text { or } \quad\left(\overline{a_{\mu \lambda}}\right)=-\left(a_{\lambda \mu}\right) . \tag{7.3}
\end{equation*}
$$

This means that $\mathrm{i} A$ is a 1 -form with values in the space $\operatorname{Herm}\left(\mathbb{C}^{r}, \mathbb{C}^{r}\right)$ of hermitian matrices. The identity $d^{2}\{s, t\}=0$ implies $\left\{D^{2} s, t\right\}+\left\{s, D^{2} t\right\}=0$, i.e. $\{\Theta(D) \wedge s, t\}+$ $\{s, \Theta(D) \wedge t\}=0$. Therefore $\Theta(D)^{\star}=-\Theta(D)$ and the curvature tensor $\Theta(D)$ is such that

$$
\mathrm{i} \Theta(D) \in \mathscr{C}_{2}^{\infty}(M, \operatorname{Herm}(E, E))
$$

(7.4) Special case. If $E$ is a hermitian line bundle $(r=1), D_{\upharpoonright \Omega}$ is a hermitian connection if and only if its connection form $A$ associated to any given orthonormal frame of $E_{\upharpoonright \Omega}$ is a 1-form with purely imaginary values.

If $\theta, \widetilde{\theta}: E_{\uparrow \Omega} \rightarrow \Omega$ are two such trivializations on a simply connected open subset $\Omega \subset M$, then $g=\widetilde{\theta} \circ \theta^{-1}=e^{\mathrm{i} \varphi}$ for some real phase function $\varphi \in \mathscr{C}^{\infty}(\Omega, \mathbb{R})$. The gauge transformation law can be written

$$
A=\widetilde{A}+\mathrm{i} d \varphi
$$

In this case, we see that $\mathrm{i} \Theta(D) \in \mathscr{C}_{2}^{\infty}(M, \mathbb{R})$.
(7.5) Remark. If $s, s^{\prime} \in \mathscr{C}^{\infty}(M, E)$ are two sections of $E$ along a smooth curve $\gamma$ : $[0,1] \longrightarrow M$, one can easily verify the formula

$$
\frac{d}{d t}\left\langle s(t), s^{\prime}(t)\right\rangle=\left\langle\frac{D s}{d t}, s^{\prime}\right\rangle+\left\langle s, \frac{D s^{\prime}}{d t}\right\rangle
$$

In particular, if $\left(e_{1}, \ldots, e_{r}\right)$ is a parallel frame of $E$ along $\gamma$ such that $\left(e_{1}(0), \ldots, e_{r}(0)\right)$ is orthonormal, then $\left(e_{1}(t), \ldots, e_{r}(t)\right)$ is orthonormal for all $t$. All parallel translation operators $T_{\gamma}$ defined in $\S 6$ are thus isometries of the fibers. It follows that $E$ has a flat hermitian connection $D$ if and only if $E$ can be defined by means of locally constant unitary transition automorphisms $g_{\alpha \beta}$, or equivalently if $E$ is isomorphic to the hermitian bundle $(\widetilde{M} \times S) / \pi_{1}(M)$ defined by a unitary representation of $\pi_{1}(M)$ in a hermitian vector space $S$. Such a bundle $E$ is said to be hermitian flat.

## $\S 8$. Vector Bundles and Locally Free Sheaves

We denote here by $\mathscr{E}$ the sheaf of germs of $\mathscr{C}^{\infty}$ complex functions on $M$. Let $F \longrightarrow M$ be a $\mathscr{C}^{\infty}$ complex vector bundle of rank $r$. We let $\mathscr{F}$ be the sheaf of germs of $\mathscr{C}^{\infty}$ sections of $F$, i.e. the sheaf whose space of sections on an open subset $U \subset M$ is $\mathscr{F}(U)=\mathscr{C}^{\infty}(U, F)$. It is clear that $\mathscr{F}$ is a $\mathscr{E}$-module. Furthermore, if $F_{\upharpoonright \Omega} \simeq \Omega \times \mathbb{C}^{r}$ is trivial, the sheaf $\mathscr{F}_{\mid \Omega}$ is isomorphic to $\mathscr{E}_{\uparrow \Omega}^{r}$ as a $\mathscr{E}_{\uparrow \Omega}$-module.
(8.1) Definition. A sheaf $\mathscr{S}$ of modules over a sheaf of rings $\mathscr{R}$ is said to be locally free of rank $k$ if every point in the base has a neighborhood $\Omega$ such that $\mathscr{S}_{\uparrow \Omega}$ is $\mathscr{R}$-isomorphic to $\mathscr{R}_{\Gamma \Omega}^{k}$.

Suppose that $\mathscr{S}$ is a locally free $\mathscr{E}$-module of rank $r$. There exists a covering $\left(V_{\alpha}\right)_{\alpha \in I}$ of $M$ and sheaf isomorphisms

$$
\theta_{\alpha}: \mathscr{S}_{\mid V_{\alpha}} \longrightarrow \mathscr{E}_{\Gamma V_{\alpha}}^{r}
$$

Then we have transition isomorphisms $g_{\alpha \beta}=\theta_{\alpha} \circ \theta_{\beta}^{-1}: \mathscr{E}^{r} \rightarrow \mathscr{E}^{r}$ defined on $V_{\alpha} \cap V_{\beta}$, and such an isomorphism is the multiplication by an invertible matrix with $\mathscr{C}^{\infty}$ coefficients on $V_{\alpha} \cap V_{\beta}$. The concepts of vector bundle and of locally free $\mathscr{E}$-module are thus completely equivalent.

Assume now that $F \longrightarrow M$ is a line bundle $(r=1)$. Then every collection of transition automorphisms $g=\left(g_{\alpha \beta}\right)$ defines a Čech 1-cocycle with values in the multiplicative sheaf $\mathscr{E}^{\star}$ of invertible $\mathscr{C}^{\infty}$ functions on $M$. In fact the definition of the Čech differential (cf. (IV-5.1)) gives $(\delta g)_{\alpha \beta \gamma}=g_{\beta \gamma} g_{\alpha \gamma}^{-1} g_{\alpha \beta}$, and we have $\delta g=1$ in view of (1.1). Let $\theta_{\alpha}^{\prime}$ be another family of trivializations and $\left(g_{\alpha \beta}^{\prime}\right)$ the associated cocycle (it is no loss of generality to assume that both are defined on the same covering since we may otherwise take a refinement). Then we have

$$
\theta_{\alpha}^{\prime} \circ \theta_{\alpha}^{-1}: V_{\alpha} \times \mathbb{C} \longrightarrow V_{\alpha} \times \mathbb{C}, \quad(x, \xi) \longmapsto\left(x, u_{\alpha}(x) \xi\right), \quad u_{\alpha} \in \mathscr{E}^{\star}\left(V_{\alpha}\right)
$$

It follows that $g_{\alpha \beta}=g_{\alpha \beta}^{\prime} u_{\alpha}^{-1} u_{\beta}$, i.e. that the Čech 1-cocycles $g, g^{\prime}$ differ only by the Čech 1 -coboundary $\delta u$. Therefore, there is a well defined map which associates to every line bundle $F$ over $M$ the Cech cohomology class $\{g\} \in H^{1}\left(M, \mathscr{E}^{\star}\right)$ of its cocycle of transition
automorphisms. It is easy to verify that the cohomology classes associated to two line bundles $F, F^{\prime}$ are equal if and only if these bundles are isomorphic: if $g=g^{\prime} \cdot \delta u$, then the collection of maps

$$
\begin{aligned}
F_{\upharpoonright V_{\alpha}} \xrightarrow{\theta_{\alpha}} V_{\alpha} \times \mathbb{C} & \longrightarrow V_{\alpha} \times \mathbb{C} \xrightarrow{\theta_{\alpha}^{\prime-1}} F_{\upharpoonright V_{\alpha}}^{\prime} \\
(x, \xi) & \longmapsto\left(x, u_{\alpha}(x) \xi\right)
\end{aligned}
$$

defines a global isomorphism $F \rightarrow F^{\prime}$. It is clear that the multiplicative group structure on $H^{1}\left(M, \mathscr{E}^{\star}\right)$ corresponds to the tensor product of line bundles (the inverse of a line bundle being given by its dual). We may summarize this discussion by the following:
(8.2) Theorem. The group of isomorphism classes of complex $\mathscr{C}^{\infty}$ line bundles is in one-to-one correspondence with the Čech cohomology group $H^{1}\left(M, \mathscr{E}^{\star}\right)$.

## § 9. First Chern Class

Throughout this section, we assume that $E$ is a complex line bundle (that is, $\mathrm{rk} E=$ $r=1)$. Let D be a connection on $E$. By (3.3), $\Theta(D)$ is a closed 2 -form on $M$. Moreover, if $D^{\prime}$ is another connection on $E$, then (2.2) shows that $D^{\prime}=D+\Gamma \wedge \bullet$ where $\Gamma \in$ $\mathscr{C}_{1}^{\infty}(M, \mathbb{C})$. By (3.3), we get

$$
\begin{equation*}
\Theta\left(D^{\prime}\right)=\Theta(D)+d \Gamma \tag{9.1}
\end{equation*}
$$

This formula shows that the De Rham class $\{\Theta(D)\} \in H_{D R}^{2}(M, \mathbb{C})$ does not depend on the particular choice of $D$. If $D$ is chosen to be hermitian with respect to a given hermitian metric on $E$ (such a connection can always be constructed by means of a partition of unity) then $\mathrm{i} \Theta(D)$ is a real 2-form, thus $\{\mathrm{i} \Theta(D)\} \in H_{D R}^{2}(M, \mathbb{R})$. Consider now the one-to-one correspondence given by Th. 8.2:

$$
\begin{aligned}
&\text { \{isomorphism classes of line bundles }\} \longrightarrow H^{1}\left(M, \mathscr{E}^{\star}\right) \\
& \text { class }\{E\} \text { defined by the cocycle }\left(g_{\alpha \beta}\right) \longmapsto \text { class of }\left(g_{\alpha \beta}\right) .
\end{aligned}
$$

Using the exponential exact sequence of sheaves

$$
\begin{aligned}
0 \longrightarrow \mathbb{Z} \longrightarrow \mathscr{E} & \longrightarrow \mathscr{E}^{\star} \longrightarrow 1 \\
f & e^{2 \pi \mathrm{i} f}
\end{aligned}
$$

and the fact that $H^{1}(M, \mathscr{E})=H^{2}(M, \mathscr{E})=0$, we obtain:
(9.2) Theorem and Definition. The coboundary morphism

$$
H^{1}\left(M, \mathscr{E}^{\star}\right) \xrightarrow{\partial} H^{2}(M, \mathbb{Z})
$$

is an isomorphism. The first Chern class of a line bundle $E$ is the image $c_{1}(E)$ in $H^{2}(M, \mathbb{Z})$ of the Čech cohomology class of the 1-cocycle $\left(g_{\alpha \beta}\right)$ associated to $E$ :

$$
\begin{equation*}
c_{1}(E)=\partial\left\{\left(g_{\alpha \beta}\right)\right\} \tag{9.3}
\end{equation*}
$$

Consider the natural morphism

$$
\begin{equation*}
H^{2}(M, \mathbb{Z}) \longrightarrow H^{2}(M, \mathbb{R}) \simeq H_{D R}^{2}(M, \mathbb{R}) \tag{9.4}
\end{equation*}
$$

where the isomorphism $\simeq$ is that given by the De Rham-Weil isomorphism theorem and the sign convention of Formula (IV-6.11).
(9.5) Theorem. The image of $c_{1}(E)$ in $H_{D R}^{2}(M, \mathbb{R})$ under (9.4) coincides with the De Rham cohomology class $\left\{\frac{\mathrm{i}}{2 \pi} \Theta(D)\right\}$ associated to any (hermitian) connection $D$ on $E$.

Proof. Choose an open covering $\left(V_{\alpha}\right)_{\alpha \in I}$ of $M$ such that $E$ is trivial on each $V_{\alpha}$, and such that all intersections $V_{\alpha} \cap V_{\beta}$ are simply connected (as in $\S$ IV-6, choose the $V_{\alpha}$ to be small balls relative to a given locally finite covering of $M$ by coordinate patches). Denote by $A_{\alpha}$ the connection forms of $D$ with respect to a family of isometric trivializations

$$
\theta_{\alpha}: E_{\uparrow V_{\alpha}} \longrightarrow V_{\alpha} \times \mathbb{C}^{r}
$$

Let $g_{\alpha \beta} \in \mathscr{E}^{\star}\left(V_{\alpha} \cap V_{\beta}\right)$ be the corresponding transition automorphisms. Then $\left|g_{\alpha \beta}\right|=1$, and as $V_{\alpha} \cap V_{\beta}$ is simply connected, we may choose real functions $u_{\alpha \beta} \in \mathscr{E}\left(V_{\alpha} \cap V_{\beta}\right)$ such that

$$
g_{\alpha \beta}=\exp \left(2 \pi \mathrm{i} u_{\alpha \beta}\right)
$$

By definition, the first Chern class $c_{1}(E)$ is the Čech 2-cocycle

$$
\begin{aligned}
c_{1}(E) & \left.=\partial\left\{\left(g_{\alpha \beta}\right)\right\}=\left\{(\delta u)_{\alpha \beta \gamma}\right)\right\} \in H^{2}(M, \mathbb{Z}) \quad \text { where } \\
(\delta u)_{\alpha \beta \gamma} & :=u_{\beta \gamma}-u_{\alpha \gamma}+u_{\alpha \beta} .
\end{aligned}
$$

Now, if $\mathscr{E}^{q}$ (resp. $\mathscr{E}^{q}$ ) denotes the sheaf of real (resp. real $d$-closed) $q$-forms on $M$, the short exact sequences

$$
\begin{aligned}
& 0 \longrightarrow \mathscr{Q}^{1} \longrightarrow \mathscr{E}^{1} \xrightarrow{d} \mathscr{E}^{2} \longrightarrow 0 \\
& 0 \longrightarrow \mathbb{R} \longrightarrow \mathscr{E}^{0} \xrightarrow{\longrightarrow} \mathscr{E}^{1} \longrightarrow 0
\end{aligned}
$$

yield isomorphisms (with the sign convention of (IV-6.11))

$$
\begin{align*}
& H_{D R}^{2}(M, \mathbb{R}):=H^{0}\left(M, \mathscr{Z}^{2}\right) / d H^{0}\left(M, \mathscr{E}^{1}\right) \xrightarrow{-\partial} H^{1}\left(M, \mathscr{E}^{1}\right),  \tag{9.6}\\
& H^{1}\left(M, \mathscr{Z}^{1}\right) \xrightarrow{\partial} H^{2}(M, \mathbb{R}) . \tag{9.7}
\end{align*}
$$

Formula 3.4 gives $A_{\beta}=A_{\alpha}+g_{\alpha \beta}^{-1} d g_{\alpha \beta}$. Since $\Theta(D)=d A_{\alpha}$ on $V_{\alpha}$, the image of $\left\{\frac{\mathrm{i}}{2 \pi} \Theta(D)\right\}$ under (9.6) is the Čech 1-cocycle with values in $\mathscr{Z}^{1}$

$$
\left\{-\frac{\mathrm{i}}{2 \pi}\left(A_{\beta}-A_{\alpha}\right)\right\}=\left\{\frac{1}{2 \pi \mathrm{i}} g_{\alpha \beta}^{-1} d g_{\alpha \beta}\right\}=\left\{d u_{\alpha \beta}\right\}
$$

and the image of this cocycle under (9.7) is the Čech 2-cocycle $\{\delta u\}$ in $H^{2}(M, \mathbb{R})$. But $\{\delta u\}$ is precisely the image of $c_{1}(E) \in H^{2}(M, \mathbb{Z})$ in $H^{2}(M, \mathbb{R})$.

Let us assume now that $M$ is oriented and that $s \in \mathscr{C}^{\infty}(M, E)$ is transverse to the zero section of $E$, i.e. that $D s \in \operatorname{Hom}(T M, E)$ is surjective at every point of the zero set $Z:=s^{-1}(0)$. Then $Z$ is an oriented 2-codimensional submanifold of $M$ (the orientation of
$Z$ is uniquely defined by those of $M$ and $E$ ). We denote by $[Z]$ the current of integration over $Z$ and by $\{[Z]\} \in H_{D R}^{2}(M, \mathbb{R})$ its cohomology class.
(9.8) Theorem. We have $\{[Z]\}=c_{1}(E)_{\mathbb{R}}$.

Proof. Consider the differential 1-form

$$
u=s^{-1} \otimes D s \in \mathscr{C}_{1}^{\infty}(M \backslash Z, \mathbb{C})
$$

Relatively to any trivialization $\theta$ of $E_{\upharpoonright \Omega}$, one has $D_{\upharpoonright \Omega} \simeq_{\theta} d+A \wedge \bullet$, thus

$$
u_{\upharpoonright \Omega}=\frac{d \sigma}{\sigma}+A \quad \text { where } \sigma=\theta(s)
$$

It follows that $u$ has locally integrable coefficients on $M$. If $d \sigma / \sigma$ is considered as a current on $\Omega$, then

$$
d\left(\frac{d \sigma}{\sigma}\right)=d\left(\sigma^{\star} \frac{d z}{z}\right)=\sigma^{\star} d\left(\frac{d z}{z}\right)=\sigma^{\star}\left(2 \pi \mathrm{i} \delta_{0}\right)=2 \pi \mathrm{i}[Z]
$$

because of the Cauchy residue formula (cf. Lemma I-2.10) and because $\sigma$ is a submersion in a neighborhood of $Z$ (cf. (I-1.19)). Now, we have $d A=\Theta(D)$ and Th. 9.8 follows from the resulting equality:

$$
\begin{equation*}
d u=2 \pi \mathrm{i}[Z]+\Theta(E) . \tag{9.9}
\end{equation*}
$$

## $\S$ 10. Connections of Type $(1,0)$ and $(0,1)$ over Complex Manifolds

Let $X$ be a complex manifold, $\operatorname{dim}_{\mathbb{C}} X=n$ and $E$ a $\mathscr{C}^{\infty}$ vector bundle of rank $r$ over $X$; here, $E$ is not assumed to be holomorphic. We denote by $\mathscr{C}_{p, q}^{\infty}(X, E)$ the space of $\mathscr{C}^{\infty}$ sections of the bundle $\Lambda^{p, q} T^{\star} X \otimes E$. We have therefore a direct sum decomposition

$$
\mathscr{C}_{l}^{\infty}(X, E)=\bigoplus_{p+q=l} \mathscr{C}_{p, q}^{\infty}(X, E)
$$

Connections of type $(1,0)$ or $(0,1)$ are operators acting on vector valued forms, which imitate the usual operators $d^{\prime}, d^{\prime \prime}$ acting on $\mathscr{C}_{p, q}^{\infty}(X, \mathbb{C})$. More precisely, a connection of type $(1,0)$ on $E$ is a differential operator $D^{\prime}$ of order 1 acting on $\mathscr{C}_{\bullet \bullet}^{\infty}(X, E)$ and satisfying the following two properties:

$$
\begin{align*}
& D^{\prime}: \mathscr{C}_{p, q}^{\infty}(X, E) \longrightarrow \mathscr{C}_{p+1, q}^{\infty}(X, E),  \tag{10.1}\\
& D^{\prime}(f \wedge s)=d^{\prime} f \wedge s+(-1)^{\operatorname{deg} f} f \wedge D^{\prime} s \tag{10.1'}
\end{align*}
$$

for any $f \in \mathscr{C}_{p_{1}, q_{1}}^{\infty}(X, \mathbb{C}), s \in \mathscr{C}_{p_{2}, q_{2}}^{\infty}(X, E)$. The definition of a connection $D^{\prime \prime}$ of type $(0,1)$ is similar. If $\theta: E_{\uparrow \Omega} \rightarrow \Omega \times \mathbb{C}^{r}$ is a $\mathscr{C}^{\infty}$ trivialization of $E_{\uparrow \Omega}$ and if $\sigma=\left(\sigma_{\lambda}\right)=\theta(s)$, then all such connections $D^{\prime}$ and $D^{\prime \prime}$ can be written

$$
\begin{align*}
D^{\prime} s & \simeq_{\theta} d^{\prime} \sigma+A^{\prime} \wedge \sigma, \\
D^{\prime \prime} s & \simeq_{\theta} d^{\prime \prime} \sigma+A^{\prime \prime} \wedge \sigma
\end{align*}
$$

where $A^{\prime} \in \mathscr{C}_{1,0}^{\infty}\left(\Omega, \operatorname{Hom}\left(\mathbb{C}^{r}, \mathbb{C}^{r}\right)\right), A^{\prime \prime} \in \mathscr{C}_{0,1}^{\infty}\left(\Omega, \operatorname{Hom}\left(\mathbb{C}^{r}, \mathbb{C}^{r}\right)\right)$ are arbitrary forms with matrix coefficients.

It is clear that $D=D^{\prime}+D^{\prime \prime}$ is then a connection in the sense of $\S 2$; conversely any connection $D$ admits a unique decomposition $D=D^{\prime}+D^{\prime \prime}$ in terms of a (1,0)-connection and a ( 0,1 )-connection.

Assume now that $E$ has a hermitian structure and that $\theta$ is an isometry. The connection $D$ is hermitian if and only if the connection form $A=A^{\prime}+A^{\prime \prime}$ satisfies $A^{\star}=-A$, and this condition is equivalent to $A^{\prime}=-\left(A^{\prime \prime}\right)^{\star}$. From this observation, we get immediately:
(10.3) Proposition. Let $D_{0}^{\prime \prime}$ be a given ( 0,1 )-connection on a hermitian bundle $\pi: E \rightarrow X$. Then there exists a unique hermitian connection $D=D^{\prime}+D^{\prime \prime}$ such that $D^{\prime \prime}=D_{0}^{\prime \prime}$.

## § 11. Holomorphic Vector Bundles

From now on, the vector bundles $E$ in which we are interested are supposed to have a holomorphic structure:
(11.1) Definition. A vector bundle $\pi: E \rightarrow X$ is said to be holomorphic if $E$ is a complex manifold, if the projection map $\pi$ is holomorphic and if there exists a covering $\left(V_{\alpha}\right)_{\alpha \in I}$ of $X$ and a family of holomorphic trivializations $\theta_{\alpha}: E_{\mid V_{\alpha}} \rightarrow V_{\alpha} \times \mathbb{C}^{r}$.

It follows that the transition matrices $g_{\alpha \beta}$ are holomorphic on $V_{\alpha} \cap V_{\beta}$. In complete analogy with the discussion of $\S 8$, we see that the concept of holomorphic vector bundle is equivalent to the concept of locally free sheaf of modules over the ring $\mathfrak{C}$ of germs of holomorphic functions on $X$. We shall denote by $\mathscr{O}(E)$ the associated sheaf of germs of holomorphic sections of $E$. In the case $r=1$, there is a one-to-one correspondence between the isomorphism classes of holomorphic line bundles and the Čech cohomology group $H^{1}\left(X, \mathscr{O}^{\star}\right)$.
(11.2) Definition. The group $H^{1}\left(X, \mathbb{O}^{\star}\right)$ of isomorphism classes of holomorphic line bundles is called the Picard group of $X$.

If $s \in \mathscr{C}_{p, q}^{\infty}(X, E)$, the components $\sigma^{\alpha}=\left(\sigma_{\lambda}^{\alpha}\right)_{1 \leqslant \lambda \leqslant r}=\theta_{\alpha}(s)$ of $s$ under $\theta_{\alpha}$ are related by

$$
\sigma^{\alpha}=g_{\alpha \beta} \cdot \sigma^{\beta} \quad \text { on } \quad V_{\alpha} \cap V_{\beta} .
$$

Since $d^{\prime \prime} g_{\alpha \beta}=0$, it follows that

$$
d^{\prime \prime} \sigma^{\alpha}=g_{\alpha \beta} \cdot d^{\prime \prime} \sigma^{\beta} \quad \text { on } \quad V_{\alpha} \cap V_{\beta} .
$$

The collection of forms $\left(d^{\prime \prime} \sigma^{\alpha}\right)$ therefore corresponds to a unique global $(p, q+1)$-form $d^{\prime \prime} s$ such that $\theta_{\alpha}\left(d^{\prime \prime} s\right)=d^{\prime \prime} \sigma^{\alpha}$, and the operator $d^{\prime \prime}$ defined in this way is a $(0,1)$-connection on $E$.
(11.3) Definition. The operator $d^{\prime \prime}$ is called the canonical ( 0,1 )-connection of the holomorphic bundle E.

It is clear that $d^{\prime 2}=0$. Therefore, for any integer $p=0,1, \ldots, n$, we get a complex

$$
\mathscr{C}_{p, 0}^{\infty}(X, E) \xrightarrow{d^{\prime \prime}} \cdots \longrightarrow \mathscr{C}_{p, q}^{\infty}(X, E) \xrightarrow{d^{\prime \prime}} \mathscr{C}_{p, q+1}^{\infty}(X, E) \longrightarrow \cdots
$$

known as the Dolbeault complex of $(p, \bullet)$-forms with values in $E$.
(11.4) Notation. The q-th cohomology group of the Dolbeault complex is denoted $H^{p, q}(X, E)$ and is called the $(p, q)$ Dolbeault cohomology group with values in $E$.

The Dolbeault-Grothendieck lemma I-2.11 shows that the complex of sheaves $d^{\prime \prime}$ : $\mathscr{C}_{0, \bullet}^{\infty}(X, E)$ is a soft resolution of the sheaf $\mathfrak{O}(E)$. By the De Rham-Weil isomorphism theorem IV-6.4, we get:
(11.5) Proposition. $H^{0, q}(X, E) \simeq H^{q}(X, \mathbb{B}(E))$.

Most often, we will identify the locally free sheaf $\mathscr{O}(E)$ and the bundle $E$ itself ; the above sheaf cohomology group will therefore be simply denoted $H^{q}(X, E)$. Another standard notation in analytic or algebraic geometry is:
(11.6) Notation. If $X$ is a complex manifold, $\Omega_{X}^{p}$ denotes the vector bundle $\Lambda^{p} T^{\star} X$ or its sheaf of sections.

It is clear that the complex $\mathscr{C}_{p, \bullet}^{\infty}(X, E)$ is identical to the complex $\mathscr{C}_{0, \bullet}^{\infty}\left(X, \Omega_{X}^{p} \otimes E\right)$, therefore we obtain a canonical isomorphism:
(11.7) Dolbeault isomorphism. $H^{p, q}(X, E) \simeq H^{q}\left(X, \Omega_{X}^{p} \otimes E\right)$.

In particular, $H^{p, 0}(X, E)$ is the space of global holomorphic sections of the bundle $\Omega_{X}^{p} \otimes E$.

## § 12. Chern Connection

Let $\pi: E \rightarrow X$ be a hermitian holomorphic vector bundle of rank $r$. By Prop. 10.3, there exists a unique hermitian connection $D$ such that $D^{\prime \prime}=d^{\prime \prime}$.
(12.1) Definition. The unique hermitian connection $D$ such that $D^{\prime \prime}=d^{\prime \prime}$ is called the Chern connection of $E$. The curvature tensor of this connection will be denoted by $\Theta(E)$ and is called the Chern curvature tensor of $E$.

Let us compute $D$ with respect to an arbitrary holomorphic trivialization $\theta: E_{\upharpoonright \Omega} \rightarrow$ $\Omega \times \mathbb{C}^{r}$. Let $H=\left(h_{\lambda \mu}\right)_{1 \leqslant \lambda, \mu \leqslant r}$ denote the hermitian matrix with $\mathscr{C}^{\infty}$ coefficients representing the metric along the fibers of $E_{\mid \Omega}$. For any $s, t \in \mathscr{C}_{\bullet \bullet}^{\infty}(X, E)$ and $\sigma=\theta(s), \tau=$ $\theta(t)$ one can write

$$
\{s, t\}=\sum_{\lambda, \mu} h_{\lambda \mu} \sigma_{\lambda} \wedge \bar{\tau}_{\mu}=\sigma^{\dagger} \wedge H \bar{\tau},
$$

where $\sigma^{\dagger}$ is the transposed matrix of $\sigma$. It follows that

$$
\begin{aligned}
\{D s, t\}+ & (-1)^{\operatorname{deg} s}\{s, D t\}=d\{s, t\} \\
& =(d \sigma)^{\dagger} \wedge H \bar{\tau}+(-1)^{\operatorname{deg} \sigma} \sigma^{\dagger} \wedge(d H \wedge \bar{\tau}+H \overline{d \tau}) \\
& =\left(d \sigma+\bar{H}^{-1} d^{\prime} \bar{H} \wedge \sigma\right)^{\dagger} \wedge H \bar{\tau}+(-1)^{\operatorname{deg} \sigma} \sigma^{\dagger} \wedge\left(\overline{\left.d \tau+\bar{H}^{-1} d^{\prime} \bar{H} \wedge \tau\right)}\right.
\end{aligned}
$$

using the fact that $d H=d^{\prime} H+\overline{d^{\prime} \bar{H}}$ and $\bar{H}^{\dagger}=H$. Therefore the Chern connection $D$ coincides with the hermitian connection defined by

$$
\begin{align*}
& D s \simeq_{\theta} d \sigma+\bar{H}^{-1} d^{\prime} \bar{H} \wedge \sigma,  \tag{12.2}\\
& D^{\prime} \simeq_{\theta} d^{\prime}+\bar{H}^{-1} d^{\prime} \bar{H} \wedge \bullet=\bar{H}^{-1} d^{\prime}(\bar{H} \bullet), \quad D^{\prime \prime}=d^{\prime \prime} \tag{12.3}
\end{align*}
$$

It is clear from this relations that $D^{\prime 2}=D^{\prime \prime 2}=0$. Consequently $D^{2}$ is given by to $D^{2}=D^{\prime} D^{\prime \prime}+D^{\prime \prime} D^{\prime}$, and the curvature tensor $\Theta(E)$ is of type $(1,1)$. Since $d^{\prime} d^{\prime \prime}+d^{\prime \prime} d^{\prime}=0$, we get

$$
\left(D^{\prime} D^{\prime \prime}+D^{\prime \prime} D^{\prime}\right) s \simeq_{\theta} \bar{H}^{-1} d^{\prime} \bar{H} \wedge d^{\prime \prime} \sigma+d^{\prime \prime}\left(\bar{H}^{-1} d^{\prime} \bar{H} \wedge \sigma\right)=d^{\prime \prime}\left(\bar{H}^{-1} d^{\prime} \bar{H}\right) \wedge \sigma .
$$

(12.4) Theorem. The Chern curvature tensor is such that

$$
\mathrm{i} \Theta(E) \in \mathscr{C}_{1,1}^{\infty}(X, \operatorname{Herm}(E, E))
$$

If $\theta: E_{\uparrow} \rightarrow \Omega \times \mathbb{C}^{r}$ is a holomorphic trivialization and if $H$ is the hermitian matrix representing the metric along the fibers of $E_{\mid \Omega}$, then

$$
\mathrm{i} \Theta(E)=\mathrm{i} d^{\prime \prime}\left(\bar{H}^{-1} d^{\prime} \bar{H}\right) \quad \text { on } \quad \Omega .
$$

Let $\left(e_{1}, \ldots, e_{r}\right)$ be a $\mathscr{C}^{\infty}$ orthonormal frame of $E$ over a coordinate patch $\Omega \subset X$ with complex coordinates $\left(z_{1}, \ldots, z_{n}\right)$. On $\Omega$ the Chern curvature tensor can be written

$$
\begin{equation*}
\mathrm{i} \Theta(E)=\mathrm{i} \sum_{1 \leqslant j, k \leqslant n, 1 \leqslant \lambda, \mu \leqslant r} c_{j k \lambda \mu} d z_{j} \wedge d \bar{z}_{k} \otimes e_{\lambda}^{\star} \otimes e_{\mu} \tag{12.5}
\end{equation*}
$$

for some coefficients $c_{j k \lambda \mu} \in \mathbb{C}$. The hermitian property of $i \Theta(E)$ means that $\bar{c}_{j k \lambda \mu}=$ $c_{k j \mu \lambda}$.
(12.6) Special case. When $r=\operatorname{rank} E=1$, the hermitian matrix $H$ is a positive function which we write $H=e^{-\varphi}, \varphi \in \mathscr{C}^{\infty}(\Omega, \mathbb{R})$. By the above formulas we get

$$
\begin{gather*}
D^{\prime} \simeq_{\theta} d^{\prime}-d^{\prime} \varphi \wedge \bullet=e^{\varphi} d^{\prime}\left(e^{-\varphi} \bullet\right),  \tag{12.7}\\
\mathrm{i} \Theta(E)=\mathrm{i} d^{\prime} d^{\prime \prime} \varphi \quad \text { on } \quad \Omega . \tag{12.8}
\end{gather*}
$$

Especially, we see that $\mathrm{i} \Theta(E)$ is a closed real $(1,1)$-form on $X$.
(12.9) Remark. In general, it is not possible to find local frames $\left(e_{1}, \ldots, e_{r}\right)$ of $E_{\Gamma \Omega}$ that are simultaneously holomorphic and orthonormal. In fact, we have in this case $H=\left(\delta_{\lambda \mu}\right)$, so a necessary condition for the existence of such a frame is that $\Theta(E)=0$ on $\Omega$. Conversely, if $\Theta(E)=0$, Prop. 6.7 and Rem. 7.5 show that $E$ possesses local orthonormal parallel frames $\left(e_{\lambda}\right)$; we have in particular $D^{\prime \prime} e_{\lambda}=0$, so $\left(e_{\lambda}\right)$ is holomorphic; such a bundle $E$ arising from a unitary representation of $\pi_{1}(X)$ is said to be hermitian flat. The next proposition shows in a more local way that the Chern curvature tensor is the obstruction to the existence of orthonormal holomorphic frames: a holomorphic frame can be made "almost orthonormal" only up to curvature terms of order 2 in a neighborhood of any point.
(12.10) Proposition. For every point $x_{0} \in X$ and every coordinate system $\left(z_{j}\right)_{1 \leqslant j \leqslant n}$ at $x_{0}$, there exists a holomorphic frame $\left(e_{\lambda}\right)_{1 \leqslant \lambda \leqslant r}$ in a neighborhood of $x_{0}$ such that

$$
\left\langle e_{\lambda}(z), e_{\mu}(z)\right\rangle=\delta_{\lambda \mu}-\sum_{1 \leqslant j, k \leqslant n} c_{j k \lambda \mu} z_{j} \bar{z}_{k}+O\left(|z|^{3}\right)
$$

where ( $c_{j k \lambda \mu}$ ) are the coefficients of the Chern curvature tensor $\Theta(E)_{x_{0}}$. Such a frame $\left(e_{\lambda}\right)$ is called a normal coordinate frame at $x_{0}$.
Proof. Let $\left(h_{\lambda}\right)$ be a holomorphic frame of $E$. After replacing $\left(h_{\lambda}\right)$ by suitable linear combinations with constant coefficients, we may assume that $\left(h_{\lambda}\left(x_{0}\right)\right)$ is an orthonormal basis of $E_{x_{0}}$. Then the inner products $\left\langle h_{\lambda}, h_{\mu}\right\rangle$ have an expansion

$$
\left\langle h_{\lambda}(z), h_{\mu}(z)\right\rangle=\delta_{\lambda \mu}+\sum_{j}\left(a_{j \lambda \mu} z_{j}+a_{j \lambda \mu}^{\prime} \bar{z}_{j}\right)+O\left(|z|^{2}\right)
$$

for some complex coefficients $a_{j \lambda \mu}, a_{j \lambda \mu}^{\prime}$ such that $a_{j \lambda \mu}^{\prime}=\bar{a}_{j \mu \lambda}$. Set first

$$
g_{\lambda}(z)=h_{\lambda}(z)-\sum_{j, \mu} a_{j \lambda \mu} z_{j} h_{\mu}(z) .
$$

Then there are coefficients $a_{j k \lambda \mu}, a_{j k \lambda \mu}^{\prime}, a_{j k \lambda \mu}^{\prime \prime}$ such that

$$
\begin{aligned}
\left\langle g_{\lambda}(z), g_{\mu}(z)\right\rangle & =\delta_{\lambda \mu}+O\left(|z|^{2}\right) \\
& =\delta_{\lambda \mu}+\sum_{j, k}\left(a_{j k \lambda \mu} z_{j} \bar{z}_{k}+a_{j k \lambda \mu}^{\prime} z_{j} z_{k}+a_{j k \lambda \mu}^{\prime \prime} \bar{z}_{j} \bar{z}_{k}\right)+O\left(|z|^{3}\right)
\end{aligned}
$$

The holomorphic frame $\left(e_{\lambda}\right)$ we are looking for is

$$
e_{\lambda}(z)=g_{\lambda}(z)-\sum_{j, k, \mu} a_{j k \lambda \mu}^{\prime} z_{j} z_{k} g_{\mu}(z)
$$

Since $a_{j k \lambda \mu}^{\prime \prime}=\bar{a}_{j k \mu \lambda}^{\prime}$, we easily find

$$
\begin{aligned}
\left\langle e_{\lambda}(z), e_{\mu}(z)\right\rangle & =\delta_{\lambda \mu}+\sum_{j, k} a_{j k \lambda \mu} z_{j} \bar{z}_{k}+O\left(|z|^{3}\right) \\
d^{\prime}\left\langle e_{\lambda}, e_{\mu}\right\rangle & =\left\{D^{\prime} e_{\lambda}, e_{\mu}\right\}=\sum_{j, k} a_{j k \lambda \mu} \bar{z}_{k} d z_{j}+O\left(|z|^{2}\right), \\
\Theta(E) \cdot e_{\lambda} & =D^{\prime \prime}\left(D^{\prime} e_{\lambda}\right)=\sum_{j, k, \mu} a_{j k \lambda \mu} d \bar{z}_{k} \wedge d z_{j} \otimes e_{\mu}+O(|z|),
\end{aligned}
$$

therefore $c_{j k \lambda \mu}=-a_{j k \lambda \mu}$.

## § 13. Lelong-Poincaré Equation and First Chern Class

Our goal here is to extend the Lelong-Poincaré equation III-2.15 to any meromorphic section of a holomorphic line bundle.
(13.1) Definition. A meromorphic section of a bundle $E \rightarrow X$ is a section s defined on an open dense subset of $X$, such that for every trivialization $\theta_{\alpha}: E_{\mid V_{\alpha}} \rightarrow V_{\alpha} \times \mathbb{C}^{r}$ the components of $\sigma^{\alpha}=\theta_{\alpha}(s)$ are meromorphic functions on $V_{\alpha}$.

Let $E$ be a hermitian line bundle, $s$ a meromorphic section which does not vanish on any component of $X$ and $\sigma=\theta(s)$ the corresponding meromorphic function in a trivialization $\theta: E_{\upharpoonright \Omega} \rightarrow \Omega \times \mathbb{C}$. The divisor of $s$ is the current on $X$ defined by div $s_{\upharpoonright \Omega}=$ $\operatorname{div} \sigma$ for all trivializing open sets $\Omega$. One can write $\operatorname{div} s=\sum m_{j} Z_{j}$, where the sets $Z_{j}$ are the irreducible components of the sets of zeroes and poles of $s$ (cf. § II-5). The Lelong-Poincaré equation (II-5.32) gives

$$
\frac{\mathrm{i}}{\pi} d^{\prime} d^{\prime \prime} \log |\sigma|=\sum m_{j}\left[Z_{j}\right],
$$

and from the equalities $|s|^{2}=|\sigma|^{2} e^{-\varphi}$ and $d^{\prime} d^{\prime \prime} \varphi=\Theta(E)$ we get

$$
\begin{equation*}
i d^{\prime} d^{\prime \prime} \log |s|^{2}=2 \pi \sum m_{j}\left[Z_{j}\right]-\mathrm{i} \Theta(E) \tag{13.2}
\end{equation*}
$$

This equality can be viewed as a complex analogue of (9.9) (except that here the hypersurfaces $Z_{j}$ are not necessarily smooth). In particular, if $s$ is a non vanishing holomorphic section of $E_{\lceil\Omega}$, we have

$$
\begin{equation*}
\mathrm{i} \Theta(E)=-\mathrm{i} d^{\prime} d^{\prime \prime} \log |s|^{2} \quad \text { on } \Omega \tag{13.3}
\end{equation*}
$$

(13.4) Theorem. Let $E \rightarrow X$ be a line bundle and let $s$ be a meromorphic section of $E$ which does not vanish identically on any component of $X$. If $\sum m_{j} Z_{j}$ is the divisor of $s$, then

$$
c_{1}(E)_{\mathbb{R}}=\left\{\sum m_{j}\left[Z_{j}\right]\right\} \in H^{2}(X, \mathbb{R})
$$

Proof. Apply Formula (13.2) and Th. 9.5, and observe that the bidimension (1, 1)-current $\mathrm{i} d^{\prime} d^{\prime \prime} \log |s|^{2}=d\left(\mathrm{i} d^{\prime \prime} \log |s|^{2}\right)$ has zero cohomology class.
(13.5) Example. If $\Delta=\sum m_{j} Z_{j}$ is an arbitrary divisor on $X$, we associate to $\Delta$ the sheaf $\mathscr{O}(\Delta)$ of germs of meromorphic functions $f$ such that $\operatorname{div}(f)+\Delta \geqslant 0$. Let $\left(V_{\alpha}\right)$ be a covering of $X$ and $u_{\alpha}$ a meromorphic function on $V_{\alpha}$ such that $\operatorname{div}\left(u_{\alpha}\right)=\Delta$ on $V_{\alpha}$. Then $\mathfrak{O}(\Delta)_{\mid V_{\alpha}}=u_{\alpha}^{-1} \mathfrak{O}$, thus $\mathfrak{O}(\Delta)$ is a locally free $\mathscr{C}$-module of rank 1 . This sheaf can be identified to the line bundle $E$ over $X$ defined by the cocycle $g_{\alpha \beta}:=u_{\alpha} / u_{\beta} \in \mathscr{O}^{\star}\left(V_{\alpha} \cap V_{\beta}\right)$. In fact, there is a sheaf isomorphism $\mathscr{O}(\Delta) \longrightarrow \mathscr{O}(E)$ defined by

$$
\mathfrak{O}(\Delta)(\Omega) \ni f \longmapsto s \in \mathbb{O}(E)(\Omega) \text { with } \theta_{\alpha}(s)=f u_{\alpha} \text { on } \Omega \cap V_{\alpha} .
$$

The constant meromorphic function $f=1$ induces a meromorphic section $s$ of $E$ such that $\operatorname{div} s=\operatorname{div} u_{\alpha}=\Delta$; in the special case when $\Delta \geqslant 0$, the section $s$ is holomorphic and its zero set $s^{-1}(0)$ is the support of $\Delta$. By Th. 13.4, we have

$$
\begin{equation*}
c_{1}(\mathbb{O}(\Delta))_{\mathbb{R}}=\{[\Delta]\} \tag{13.6}
\end{equation*}
$$

Let us consider the exact sequence $1 \rightarrow \mathscr{O}^{\star} \rightarrow \mathcal{M}^{\star} \rightarrow$ Div $\rightarrow 0$ already described in (II-5.36). There is a corresponding cohomology exact sequence

$$
\begin{equation*}
\mathscr{M}^{\star}(X) \longrightarrow \operatorname{Div}(X) \xrightarrow{\partial^{0}} H^{1}\left(X, \mathscr{O}^{\star}\right) \tag{13.7}
\end{equation*}
$$

The connecting homomorphism $\partial^{0}$ is equal to the map

$$
\Delta \longmapsto \text { isomorphism class of } \mathfrak{O}(\Delta)
$$

defined above. The kernel of this map consists of divisors which are divisors of global meromorphic functions in $\mathscr{M}^{\star}(X)$. In particular, two divisors $\Delta_{1}$ and $\Delta_{2}$ give rise to isomorphic line bundles $\mathfrak{G}\left(\Delta_{1}\right) \simeq \mathscr{O}\left(\Delta_{2}\right)$ if and only if $\Delta_{2}-\Delta_{1}=\operatorname{div}(f)$ for some global meromorphic function $f \in \mathcal{M}^{\star}(X)$; such divisors are called linearly equivalent. The image of $\partial^{0}$ consists of classes of line bundles $E$ such that $E$ has a global meromorphic section which does not vanish on any component of $X$. Indeed, if $s$ is such a section and $\Delta=\operatorname{div} s$, there is an isomorphism

$$
\begin{equation*}
\mathscr{O}(\Delta) \xrightarrow{\simeq} \mathscr{O}(E), \quad f \longmapsto f s . \tag{13.8}
\end{equation*}
$$

The last result of this section is a characterization of 2 -forms on $X$ which can be written as the curvature form of a hermitian holomorphic line bundle.
(13.9) Theorem. Let $X$ be an arbitrary complex manifold.
a) For any hermitian line bundle $E$ over $M$, the Chern curvature form $\frac{i}{2 \pi} \Theta(E)$ is a closed real $(1,1)$-form whose De Rham cohomology class is the image of an integral class.
b) Conversely, let $\omega$ be a $\mathscr{C}^{\infty}$ closed real $(1,1)$-form such that $\{\omega\} \in H_{D R}^{2}(X, \mathbb{R})$ is the image of an integral class. Then there exists a hermitian line bundle $E \rightarrow X$ such that $\frac{\mathrm{i}}{2 \pi} \Theta(E)=\omega$.

Proof. a) is an immediate consequence of Formula (12.9) and Th. 9.5, so we have only to prove the converse part b). By Prop. III-1.20, there exist an open covering $\left(V_{\alpha}\right)$ of $X$ and functions $\varphi_{\alpha} \in \mathscr{C}^{\infty}\left(V_{\alpha}, \mathbb{R}\right)$ such that $\frac{\mathrm{i}}{2 \pi} d^{\prime} d^{\prime \prime} \varphi_{\alpha}=\omega$ on $V_{\alpha}$. It follows that the function $\varphi_{\beta}-\varphi_{\alpha}$ is pluriharmonic on $V_{\alpha} \cap V_{\beta}$. If $\left(V_{\alpha}\right)$ is chosen such that the intersections $V_{\alpha} \cap V_{\beta}$ are simply connected, then Th. I-3.35 yields holomorphic functions $f_{\alpha \beta}$ on $V_{\alpha} \cap V_{\beta}$ such that

$$
2 \operatorname{Re} f_{\alpha \beta}=\varphi_{\beta}-\varphi_{\alpha} \quad \text { on } \quad V_{\alpha} \cap V_{\beta} .
$$

Now, our aim is to prove (roughly speaking) that $\left(\exp \left(-f_{\alpha \beta}\right)\right)$ is a cocycle in $\mathscr{O}^{\star}$ that defines the line bundle $E$ we are looking for. The Čech differential $(\delta f)_{\alpha \beta \gamma}=f_{\beta \gamma}-f_{\alpha \gamma}+$ $f_{\alpha \beta}$ takes values in the constant sheaf $\mathbb{i} \mathbb{R}$ because

$$
2 \operatorname{Re}(\delta f)_{\alpha \beta \gamma}=\left(\varphi_{\gamma}-\varphi_{\beta}\right)-\left(\varphi_{\gamma}-\varphi_{\alpha}\right)+\left(\varphi_{\beta}-\varphi_{\alpha}\right)=0
$$

Consider the real 1-forms $A_{\alpha}=\frac{\mathrm{i}}{4 \pi}\left(d^{\prime \prime} \varphi_{\alpha}-d^{\prime} \varphi_{\alpha}\right)$. As $d^{\prime}\left(\varphi_{\beta}-\varphi_{\alpha}\right)$ is equal to $d^{\prime}\left(f_{\alpha \beta}+\right.$ $\left.\bar{f}_{\alpha \beta}\right)=d f_{\alpha \beta}$, we get

$$
(\delta A)_{\alpha \beta}=A_{\beta}-A_{\alpha}=\frac{\mathrm{i}}{4 \pi} d\left(\bar{f}_{\alpha \beta}-f_{\alpha \beta}\right)=\frac{1}{2 \pi} d \operatorname{Im} f_{\alpha \beta} .
$$

Since $\omega=d A_{\alpha}$, it follows by (9.6) and (9.7) that the Čech cohomology class $\left\{\delta\left(\frac{1}{2 \pi} \operatorname{Im} f_{\alpha \beta}\right)\right\}$ is equal to $\{\omega\} \in H^{2}(X, \mathbb{R})$, which is by hypothesis the image of a 2-cocycle $\left(n_{\alpha \beta \gamma}\right) \in$ $H^{2}(X, \mathbb{Z})$. Thus we can write

$$
\delta\left(\frac{1}{2 \pi} \operatorname{Im} f_{\alpha \beta}\right)=\left(n_{\alpha \beta \gamma}\right)+\delta\left(c_{\alpha \beta}\right)
$$

for some 1-chain ( $c_{\alpha \beta}$ ) with values in $\mathbb{R}$. If we replace $f_{\alpha \beta}$ by $f_{\alpha \beta}-2 \pi \mathrm{i} c_{\alpha \beta}$, then we can achieve $c_{\alpha \beta}=0$, so $\delta\left(f_{\alpha \beta}\right) \in 2 \pi \mathrm{i} \mathbb{Z}$ and $g_{\alpha \beta}:=\exp \left(-f_{\alpha \beta}\right)$ will be a cocycle with values in $\mathfrak{O}^{\star}$. Since

$$
\varphi_{\beta}-\varphi_{\alpha}=2 \operatorname{Re} f_{\alpha \beta}=-\log \left|g_{\alpha \beta}\right|^{2}
$$

the line bundle $E$ associated to this cocycle admits a global hermitian metric defined in every trivialization by the matrix $H_{\alpha}=\left(\exp \left(-\varphi_{\alpha}\right)\right)$ and therefore

$$
\frac{\mathrm{i}}{2 \pi} \Theta(E)=\frac{\mathrm{i}}{2 \pi} d^{\prime} d^{\prime \prime} \varphi_{\alpha}=\omega \quad \text { on } \quad V_{\alpha} .
$$

## §14. Exact Sequences of Hermitian Vector Bundles

Let us consider an exact sequence of holomorphic vector bundles over $X$ :

$$
\begin{equation*}
0 \longrightarrow S \xrightarrow{j} E \xrightarrow{g} Q \longrightarrow 0 \tag{14.1}
\end{equation*}
$$

Then $E$ is said to be an extension of $S$ by $Q$. A (holomorphic, resp. $\mathscr{C}^{\infty}$ ) splitting of the exact sequence is a (holomorphic, resp. $\mathscr{C}^{\infty}$ ) homomorphism $h: Q \longrightarrow E$ which is a right inverse of the projection $E \longrightarrow Q$, i.e. such that $g \circ h=\operatorname{Id}_{Q}$.

Assume that a $\mathscr{C}^{\infty}$ hermitian metric on $E$ is given. Then $S$ and $Q$ can be endowed with the induced and quotient metrics respectively. Let us denote by $D_{E}, D_{S}, D_{Q}$ the corresponding Chern connections. The adjoint homomorphisms

$$
j^{\star}: E \longrightarrow S, \quad g^{\star}: Q \longrightarrow E
$$

are $\mathscr{C}^{\infty}$ and can be described respectively as the orthogonal projection of $E$ onto $S$ and as the orthogonal splitting of the exact sequence (14.1). They yield a $\mathscr{C}^{\infty}$ (in general non analytic) isomorphism

$$
\begin{equation*}
j^{\star} \oplus g: E \xrightarrow{\simeq} S \oplus Q \tag{14.2}
\end{equation*}
$$

(14.3) Theorem. According to the $\mathscr{C}^{\infty}$ isomorphism $j^{\star} \oplus g, D_{E}$ can be written

$$
D_{E}=\left(\begin{array}{cc}
D_{S} & -\beta^{\star} \\
\beta & D_{Q}
\end{array}\right)
$$

where $\beta \in \mathscr{C}_{1,0}^{\infty}(X, \operatorname{Hom}(S, Q))$ is called the second fundamental of $S$ in $E$ and where $\beta^{\star} \in \mathscr{C}_{0,1}^{\infty}(X, \operatorname{Hom}(Q, S))$ is the adjoint of $\beta$. Furthermore, the following identities hold:
a) $\quad D_{\operatorname{Hom}(S, E)}^{\prime} j=g^{\star} \circ \beta, \quad d^{\prime \prime} j=0$;
b) $D_{\operatorname{Hom}(E, Q)}^{\prime} g=-\beta \circ j^{\star}, d^{\prime \prime} g=0$;
c) $\quad D_{\operatorname{Hom}(E, S)}^{\prime} j^{\star}=0, \quad d^{\prime \prime} j^{\star}=\beta^{\star} \circ g$;
d) $D_{\operatorname{Hom}(Q, E)}^{\prime} g^{\star}=0, \quad d^{\prime \prime} g^{\star}=-j \circ \beta^{\star}$.

Proof. If we define $\nabla_{E} \simeq D_{S} \oplus D_{Q}$ via (14.2), then $\nabla_{E}$ is a hermitian connection on $E$. By (7.3), we have therefore $D_{E}=\nabla_{E}+\Gamma \wedge \bullet$, where $\Gamma \in \mathscr{C}_{1}^{\infty}(X, \operatorname{Hom}(E, E))$ and $\Gamma^{\star}=-\Gamma$. Let us write

$$
\begin{gather*}
\Gamma=\left(\begin{array}{ll}
\alpha & \gamma \\
\beta & \delta
\end{array}\right), \quad \alpha^{\star}=-\alpha, \delta^{\star}=-\delta, \gamma=-\beta^{\star}, \\
D_{E}=\left(\begin{array}{cc}
D_{S}+\alpha & \gamma \\
\beta & D_{Q}+\delta
\end{array}\right) . \tag{14.4}
\end{gather*}
$$

For any section $u \in \mathscr{C}_{\bullet, \bullet}^{\infty}(X, E)$ we have

$$
\begin{aligned}
D_{E} u & =D_{E}\left(j j^{\star} u+g^{\star} g u\right) \\
& =j D_{S}\left(j^{\star} u\right)+g^{\star} D_{Q}(g u)+\left(D_{\operatorname{Hom}(S, E)} j\right) \wedge j^{\star} u+\left(D_{\operatorname{Hom}(E, Q)} g^{\star}\right) \wedge g u .
\end{aligned}
$$

A comparison with (14.4) yields

$$
\begin{aligned}
D_{\mathrm{Hom}(S, E)} j & =j \circ \alpha+g^{\star} \circ \beta, \\
D_{\mathrm{Hom}(E, Q)} g^{\star} & =j \circ \gamma+g^{\star} \circ \delta,
\end{aligned}
$$

Since $j$ is holomorphic, we have $d^{\prime \prime} j=j \circ \alpha^{0,1}+g^{\star} \circ \beta^{0,1}=0$, thus $\alpha^{0,1}=\beta^{0,1}=0$. But $\alpha^{\star}=-\alpha$, hence $\alpha=0$ and $\beta \in \mathscr{C}_{1,0}^{\infty}(\operatorname{Hom}(S, Q))$; identity a) follows. Similarly, we get

$$
\begin{aligned}
D_{S}\left(j^{\star} u\right) & =j^{\star} D_{E} u+\left(D_{\operatorname{Hom}(E, S)} j^{\star}\right) \wedge u, \\
D_{Q}(g u) & =g D_{E} u+\left(D_{\operatorname{Hom}(E, Q)} g\right) \wedge u,
\end{aligned}
$$

and comparison with (14.4) yields

$$
\begin{aligned}
D_{\mathrm{Hom}(E, S)} j^{\star} & =-\alpha \circ j^{\star}-\gamma \circ g=\beta^{\star} \circ g, \\
D_{\mathrm{Hom}(E, Q)} g & =-\beta \circ j^{\star}-\delta \circ g .
\end{aligned}
$$

Since $d^{\prime \prime} g=0$, we get $\delta^{0,1}=0$, hence $\delta=0$. Identities b), c), d) follow from the above computations.
(14.5) Theorem. We have $d^{\prime \prime}\left(\beta^{\star}\right)=0$, and the Chern curvature of $E$ is

$$
\Theta(E)=\left(\begin{array}{cc}
\Theta(S)-\beta^{\star} \wedge \beta & D_{\operatorname{Hom}(Q, S)}^{\prime} \beta^{\star} \\
d^{\prime \prime} \beta & \Theta(Q)-\beta \wedge \beta^{\star}
\end{array}\right) .
$$

Proof. A computation of $D_{E}^{2}$ yields

$$
D_{E}^{2}=\left(\begin{array}{cc}
D_{S}^{2}-\beta^{\star} \wedge \beta & -\left(D_{S} \circ \beta^{\star}+\beta^{\star} \circ D_{Q}\right) \\
\beta \circ D_{S}+D_{Q} \circ \beta & D_{Q}^{2}-\beta \wedge \beta^{\star}
\end{array}\right) .
$$

Formula (13.4) implies

$$
\begin{aligned}
D_{\operatorname{Hom}(S, Q)} \beta & =\beta \circ D_{S}+D_{Q} \circ \beta, \\
D_{\operatorname{Hom}(Q, S)} \beta^{\star} & =D_{S} \circ \beta^{\star}+\beta^{\star} \circ D_{Q} .
\end{aligned}
$$

Since $D_{E}^{2}$ is of type (1,1), it follows that $d^{\prime \prime} \beta^{\star}=D_{\operatorname{Hom}(Q, S)}^{\prime \prime} \beta^{\star}=0$. The proof is achieved.

A consequence of Th. 14.5 is that $\Theta(S)$ and $\Theta(Q)$ are given in terms of $\Theta(E)$ by the following formulas, where $\Theta(E)_{\mid S}, \Theta(E)_{\mid Q}$ denote the blocks in the matrix of $\Theta(E)$ corresponding to $\operatorname{Hom}(S, S)$ and $\operatorname{Hom}(Q, Q)$ :

$$
\begin{align*}
\Theta(S) & =\Theta(E)_{\upharpoonright S}+\beta^{\star} \wedge \beta  \tag{14.6}\\
\Theta(Q) & =\Theta(E)_{\upharpoonright Q}+\beta \wedge \beta^{\star} . \tag{14.7}
\end{align*}
$$

By 14.3 c ) the second fundamental form $\beta$ vanishes identically if and only if the orthogonal splitting $E \simeq S \oplus Q$ is holomorphic ; then we have $\Theta(E)=\Theta(S) \oplus \Theta(Q)$.

Next, we show that the $d^{\prime \prime}$-cohomology class $\left\{\beta^{\star}\right\} \in H^{0,1}(X, \operatorname{Hom}(Q, S))$ characterizes the isomorphism class of $E$ among all extensions of $S$ by $Q$. Two extensions $E$ and $F$ are said to be isomorphic if there is a commutative diagram of holomorphic maps

in which the rows are exact sequences. The central vertical arrow is then necessarily an isomorphism. It is easily seen that $0 \rightarrow S \rightarrow E \rightarrow Q \rightarrow 0$ has a holomorphic splitting if and only if $E$ is isomorphic to the trivial extension $S \oplus Q$.
(14.9) Proposition. The correspondence

$$
\{E\} \longmapsto\left\{\beta^{\star}\right\}
$$

induces a bijection from the set of isomorphism classes of extensions of $S$ by $Q$ onto the cohomology group $H^{1}(X, \operatorname{Hom}(Q, S))$. In particular $\left\{\beta^{\star}\right\}$ vanishes if and only if the exact sequence

$$
0 \longrightarrow S \xrightarrow{j} E \xrightarrow{g} Q \longrightarrow 0
$$

splits holomorphically.
Proof. a) The map is well defined, i.e. $\left\{\beta^{\star}\right\}$ does not depend on the choice of the hermitian metric on $E$. Indeed, a new hermitian metric produces a new $\mathscr{C}^{\infty}$ splitting $\widehat{g}^{\star}$ and a new form $\widehat{\beta}^{\star}$ such that $d^{\prime \prime} \widehat{g}^{\star}=-j \circ \widehat{\beta}^{\star}$. Then $g g^{\star}=g \widehat{g}^{\star}=\operatorname{Id}_{Q}$, thus $\widehat{g}-g=j \circ v$ for some section $v \in \mathscr{C}^{\infty}(X, \operatorname{Hom}(Q, S))$. It follows that $\widehat{\beta}^{\star}-\beta^{\star}=-d^{\prime \prime} v$. Moreover, it is clear that an isomorphic extension $F$ has the same associated form $\beta^{\star}$ if $F$ is endowed with the image of the hermitian metric of $E$.
b) The map is injective. Let $E$ and $F$ be extensions of $S$ by $Q$. Select $\mathscr{C}^{\infty}$ splittings $E, F \simeq S \oplus Q$. We endow $S, Q$ with arbitrary hermitian metrics and $E, F$ with the direct sum metric. Then we have corresponding ( 0,1 )-connections

$$
D_{E}^{\prime \prime}=\left(\begin{array}{cc}
D_{S}^{\prime \prime} & -\beta^{\star} \\
0 & D_{Q}^{\prime \prime}
\end{array}\right), \quad D_{F}^{\prime \prime}=\left(\begin{array}{cc}
D_{S}^{\prime \prime} & -\widetilde{\beta}^{\star} \\
0 & D_{Q}^{\prime \prime}
\end{array}\right)
$$

Assume that $\widetilde{\beta}^{\star}=\beta^{\star}+d^{\prime \prime} v$ for some $v \in \mathscr{C}^{\infty}(X, \operatorname{Hom}(Q, S))$. The isomorphism $\Psi$ : $E \longrightarrow F$ of class $\mathscr{C}^{\infty}$ defined by the matrix

$$
\left(\begin{array}{cc}
\operatorname{Id}_{S} & v \\
0 & \operatorname{Id}_{Q}
\end{array}\right) .
$$

is then holomorphic, because the relation $D_{S}^{\prime \prime} \circ v-v \circ D_{Q}^{\prime \prime}=d^{\prime \prime} v=\widetilde{\beta}^{\star}-\beta^{\star}$ implies

$$
\begin{aligned}
D_{\operatorname{Hom}(E, F)}^{\prime \prime} \Psi & =D_{F}^{\prime \prime} \circ \Psi-\Psi \circ D_{E}^{\prime \prime} \\
& =\left(\begin{array}{cc}
D_{S}^{\prime \prime} & -\widetilde{\beta}^{\star} \\
0 & D_{Q}^{\prime \prime}
\end{array}\right)\left(\begin{array}{cc}
\operatorname{Id}_{S} & v \\
0 & \operatorname{Id}_{Q}
\end{array}\right)-\left(\begin{array}{cc}
\operatorname{Id}_{S} & v \\
0 & \operatorname{Id}_{Q}
\end{array}\right)\left(\begin{array}{cc}
D_{S}^{\prime \prime} & -\beta^{\star} \\
0 & D_{Q}^{\prime \prime}
\end{array}\right) \\
& =\left(\begin{array}{cc}
0 & -\widetilde{\beta}^{\star}+\beta^{\star}+\left(D_{S}^{\prime \prime} \circ v-v \circ D_{Q}^{\prime \prime}\right) \\
0 & 0
\end{array}\right)=0 .
\end{aligned}
$$

Hence the extensions $E$ and $F$ are isomorphic.
c) The map is surjective. Let $\gamma$ be an arbitrary $d^{\prime \prime}$-closed ( 0,1 )-form on $X$ with values in $\operatorname{Hom}(Q, S)$. We define $E$ as the $\mathscr{C}^{\infty}$ hermitian vector bundle $S \oplus Q$ endowed with the $(0,1)$-connection

$$
D_{E}^{\prime \prime}=\left(\begin{array}{cc}
D_{S}^{\prime \prime} & \gamma \\
0 & D_{Q}^{\prime \prime}
\end{array}\right)
$$

We only have to show that this connection is induced by a holomorphic structure on $E$; then we will have $\beta^{\star}=-\gamma$. However, the Dolbeault-Grothendieck lemma implies that there is a covering of $X$ by open sets $U_{\alpha}$ on which $\gamma=d^{\prime \prime} v_{\alpha}$ for some $v_{\alpha} \in \mathscr{C}^{\infty}\left(U_{\alpha}, \operatorname{Hom}(Q, S)\right)$. Part b) above shows that the matrix

$$
\left(\begin{array}{cc}
\mathrm{Id}_{S} & v_{\alpha} \\
0 & \operatorname{Id}_{Q}
\end{array}\right)
$$

defines an isomorphism $\psi_{\alpha}$ from $E_{\left\lceil U_{\alpha}\right.}$ onto the trivial extension $(S \oplus Q)_{\mid U_{\alpha}}$ such that $D_{\operatorname{Hom}(E, S \oplus Q)}^{\prime \prime} \psi_{\alpha}=0$. The required holomorphic structure on $E_{\uparrow U_{\alpha}}$ is the inverse image of the holomorphic structure of $(S \oplus Q)_{\mid U_{\alpha}}$ by $\psi_{\alpha}$; it is independent of $\alpha$ because $v_{\alpha}-v_{\beta}$ and $\psi_{\alpha} \circ \psi_{\beta}^{-1}$ are holomorphic on $U_{\alpha} \cap U_{\beta}$.
(14.10) Remark. If $E$ and $F$ are extensions of $S$ by $Q$ such that the corresponding forms $\beta^{\star}$ and $\widetilde{\beta}^{\star}=u \circ \beta^{\star} \circ v^{-1}$ differ by $u \in H^{0}(X, \operatorname{Aut}(S)), v \in H^{0}(X, \operatorname{Aut}(Q))$, it is easy to see that the bundles $E$ and $F$ are isomorphic. To see this, we need only replace the vertical arrows representing the identity maps of $S$ and $Q$ in (14.8) by $u$ and $v$ respectively. Thus, if we want to classify isomorphism classes of bundles $E$ which are extensions of $S$ by $Q$ rather than the extensions themselves, the
set of classes is the quotient of $H^{1}(X, \operatorname{Hom}(Q, S))$ by the action of $H^{0}(X, \operatorname{Aut}(S)) \times$ $H^{0}(X, \operatorname{Aut}(Q))$. In particular, if $S, Q$ are line bundles and if $X$ is compact connected, then $H^{0}(X, \operatorname{Aut}(S)), H^{0}(X, \operatorname{Aut}(Q))$ are equal to $\mathbb{C}^{\star}$ and the set of classes is the projective space $P\left(H^{1}(X, \operatorname{Hom}(Q, S))\right)$.

## §15. Line Bundles $\mathbb{C}(k)$ over $\mathbb{P}^{n}$

## $\S$ 15.A. Algebraic properties of $\mathfrak{O}(k)$

Let $V$ be a complex vector space of dimension $n+1, n \geqslant 1$. The quotient topological space $P(V)=(V \backslash\{0\}) / \mathbb{C}^{\star}$ is called the projective space of $V$, and can be considered as the set of lines in $V$ if $\{0\}$ is added to each class $\mathbb{C}^{\star} \cdot x$. Let

$$
\begin{aligned}
\pi: V \backslash\{0\} & \longrightarrow P(V) \\
x & \longmapsto[x]=\mathbb{C}^{\star} \cdot x
\end{aligned}
$$

be the canonical projection. When $V=\mathbb{C}^{n+1}$, we simply denote $P(V)=\mathbb{P}^{n}$. The space $\mathbb{P}^{n}$ is the quotient $S^{2 n+1} / S^{1}$ of the unit sphere $S^{2 n+1} \subset \mathbb{C}^{n+1}$ by the multiplicative action of the unit circle $S^{1} \subset \mathbb{C}$, so $\mathbb{P}^{n}$ is compact. Let $\left(e_{0}, \ldots, e_{n}\right)$ be a basis of $V$, and let $\left(x_{0}, \ldots, x_{n}\right)$ be the coordinates of a vector $x \in V \backslash\{0\}$. Then $\left(x_{0}, \ldots, x_{n}\right)$ are called the homogeneous coordinates of $[x] \in P(V)$. The space $P(V)$ can be covered by the open sets $\Omega_{j}$ defined by $\Omega_{j}=\left\{[x] \in P(V) ; x_{j} \neq 0\right\}$ and there are homeomorphisms

$$
\begin{aligned}
\tau_{j}: & \Omega_{j} \longrightarrow \mathbb{C}^{n} \\
& {[x] \longmapsto\left(z_{0}, \ldots, \widehat{z_{j}}, \ldots, z_{n}\right), \quad z_{l}=x_{l} / x_{j} \text { for } l \neq j }
\end{aligned}
$$

The collection $\left(\tau_{j}\right)$ defines a holomorphic atlas on $P(V)$, thus $P(V)=\mathbb{P}^{n}$ is a compact $n$-dimensional complex analytic manifold.

Let $\underline{V}$ be the trivial bundle $P(V) \times V$. We denote by $\mathscr{O}(-1) \subset \underline{V}$ the tautological line subbundle

$$
\begin{equation*}
\mathfrak{O}(-1)=\{([x], \xi) \in P(V) \times V ; \xi \in \mathbb{C} \cdot x\} \tag{15.1}
\end{equation*}
$$

such that $\mathscr{O}(-1)_{[x]}=\mathbb{C} \cdot x \subset V, x \in V \backslash\{0\}$. Then $\mathscr{O}(-1)_{\mid \Omega_{j}}$ admits a non vanishing holomorphic section

$$
[x] \longrightarrow \varepsilon_{j}([x])=x / x_{j}=z_{0} e_{0}+\ldots+e_{j}+z_{j+1} e_{j+1}+\ldots+z_{n} e_{n}
$$

and this shows in particular that $\mathbb{O}(-1)$ is a holomorphic line bundle.
(15.2) Definition. For every $k \in \mathbb{Z}$, the line bundle $\mathscr{O}(k)$ is defined by

$$
\begin{aligned}
\mathscr{O}(1) & =\mathfrak{O}(-1)^{\star}, \quad \mathscr{O}(0)=P(V) \times \mathbb{C}, \\
\mathfrak{O}(k) & =\mathscr{O}(1)^{\otimes k}=\mathfrak{O}(1) \otimes \cdots \otimes \mathbb{O}(1) \quad \text { for } \quad k \geqslant 1, \\
\mathscr{O}(-k) & =\mathscr{O}(-1)^{\otimes k} \quad \text { for } \quad k \geqslant 1
\end{aligned}
$$

We also introduce the quotient vector bundle $H=\underline{V} / \mathscr{O}(-1)$ of rank $n$. Therefore we have canonical exact sequences of vector bundles over $P(V)$ :

$$
\begin{equation*}
0 \rightarrow \mathbb{O}(-1) \rightarrow \underline{V} \rightarrow H \rightarrow 0, \quad 0 \rightarrow H^{\star} \rightarrow \underline{V}^{\star} \rightarrow \mathbb{O}(1) \rightarrow 0 . \tag{15.3}
\end{equation*}
$$

The total manifold of the line bundle $\mathfrak{O}(-1)$ gives rise to the so called monoidal transformation, or Hopf $\sigma$-process:
(15.4) Lemma. The holomorphic map $\mu: \mathfrak{O}(-1) \rightarrow V$ defined by

$$
\mu: \overparen{O}(-1) \hookrightarrow \underline{V}=P(V) \times V \xrightarrow{p r_{2}} V
$$

sends the zero section $P(V) \times\{0\}$ of $\mathbb{O}(-1)$ to the point $\{0\}$ and induces a biholomorphism of $\mathscr{O}(-1) \backslash(P(V) \times\{0\})$ onto $V \backslash\{0\}$.

Proof. The inverse map $\mu^{-1}: V \backslash\{0\} \longrightarrow \mathscr{O}(-1)$ is clearly defined by

$$
\mu^{-1}: x \longmapsto([x], x) .
$$

The space $H^{0}\left(\mathbb{P}^{n}, \mathscr{O}(k)\right)$ of global holomorphic sections of $\mathfrak{O}(k)$ can be easily computed by means of the above map $\mu$.
(15.5) Theorem. $H^{0}(P(V), \mathscr{O}(k))=0$ for $k<0$, and there is a canonical isomorphism

$$
H^{0}(P(V), \overparen{O}(k)) \simeq S^{k} V^{\star}, \quad k \geqslant 0
$$

where $S^{k} V^{\star}$ denotes the $k$-th symmetric power of $V^{\star}$.
(15.6) Corollary. We have $\operatorname{dim} H^{0}\left(\mathbb{P}^{n}, \mathscr{O}(k)\right)=\binom{n+k}{n}$ for $k \geqslant 0$, and this group is 0 for $k<0$.

Proof. Assume first that $k \geqslant 0$. There exists a canonical morphism

$$
\Phi: S^{k} V^{\star} \longrightarrow H^{0}(P(V), \mathscr{O}(k)) ;
$$

indeed, any element $a \in S^{k} V^{\star}$ defines a homogeneous polynomial of degree $k$ on $V$ and thus by restriction to $\mathscr{O}(-1) \subset \underline{V}$ a section $\Phi(a)=\widetilde{a}$ of $\left(\mathscr{O}(-1)^{\star}\right)^{\otimes k}=\mathscr{O}(k)$; in other words $\Phi$ is induced by the $k$-th symmetric power $S^{k} \underline{V}^{\star} \rightarrow \mathscr{O}(k)$ of the canonical morphism $\underline{V}^{\star} \rightarrow \mathscr{O}(1)$ in (15.3).

Assume now that $k \in \mathbb{Z}$ is arbitrary and that $s$ is a holomorphic section of $\mathscr{O}(k)$. For every $x \in V \backslash\{0\}$ we have $s([x]) \in \mathscr{O}(k)_{[x]}$ and $\mu^{-1}(x) \in \mathscr{O}(-1)_{[x]}$. We can therefore associate to $s$ a holomorphic function on $V \backslash\{0\}$ defined by

$$
f(x)=s([x]) \cdot \mu^{-1}(x)^{k}, \quad x \in V \backslash\{0\} .
$$

Since $\operatorname{dim} V=n+1 \geqslant 2, f$ can be extended to a holomorphic function on $V$ and $f$ is clearly homogeneous of degree $k$ ( $\mu$ and $\mu^{-1}$ are homogeneous of degree 1). It follows
that $f=0, s=0$ if $k<0$ and that $f$ is a homogeneous polynomial of degree $k$ on $V$ if $k \geqslant 0$. Thus, there exists a unique element $a \in S^{k} V^{\star}$ such that

$$
f(x)=a \cdot x^{k}=\widetilde{a}([x]) \cdot \mu^{-1}(x)^{k}
$$

Therefore $\Phi$ is an isomorphism.
The tangent bundle on $\mathbb{P}^{n}$ is closely related to the bundles $H$ and $\mathbb{O}(1)$ as shown by the following proposition.
(15.7) Proposition. There is a canonical isomorphism of bundles

$$
T P(V) \simeq H \otimes \mathscr{O}(1)
$$

Proof. The differential $d \pi_{x}$ of the projection $\pi: V \backslash\{0\} \rightarrow P(V)$ may be considered as a map

$$
d \pi_{x}: V \rightarrow T_{[x]} P(V)
$$

As $d \pi_{x}(x)=0, d \pi_{x}$ can be factorized through $V / \mathbb{C} \cdot x=V / \mathscr{O}(-1)_{[x]}=H_{[x]}$. Hence we get an isomorphism

$$
d \widetilde{\pi}_{x}: H_{[x]} \longrightarrow T_{[x]} P(V),
$$

but this isomorphism depends on $x$ and not only on the base point $[x]$ in $P(V)$. The formula $\pi(\lambda x+\xi)=\pi\left(x+\lambda^{-1} \xi\right), \lambda \in \mathbb{C}^{\star}, \xi \in V$, shows that $d \pi_{\lambda x}=\lambda^{-1} d \pi_{x}$, hence the map

$$
d \widetilde{\pi}_{x} \otimes \mu^{-1}(x): \quad H_{[x]} \longrightarrow(T P(V) \otimes \mathscr{O}(-1))_{[x]}
$$

depends only on $[x]$. Therefore $H \simeq T P(V) \otimes \mathbb{O}(-1)$.

## $\S$ 15.B. Curvature of the Tautological Line Bundle

Assume now that $V$ is a hermitian vector space. Then (15.3) yields exact sequences of hermitian vector bundles. We shall compute the curvature of $\mathbb{O}(1)$ and $H$.

Let $a \in P(V)$ be fixed. Choose an orthonormal basis $\left(e_{0}, e_{1}, \ldots, e_{n}\right)$ of $V$ such that $a=\left[e_{0}\right]$. Consider the embedding

$$
\mathbb{C}^{n} \longleftrightarrow P(V), \quad 0 \longmapsto a
$$

which sends $z=\left(z_{1}, \ldots, z_{n}\right)$ to $\left[e_{0}+z_{1} e_{1}+\cdots+z_{n} e_{n}\right]$. Then

$$
\varepsilon(z)=e_{0}+z_{1} e_{1}+\cdots+z_{n} e_{n}
$$

defines a non-zero holomorphic section of $\mathscr{O}(-1)_{\mid \mathbb{C}^{n}}$ and Formula (13.3) for $\Theta(\mathscr{O}(1))=$ $-\Theta(\mathbb{O}(-1))$ implies

$$
\begin{align*}
\Theta(\mathscr{O}(1)) & =d^{\prime} d^{\prime \prime} \log |\varepsilon(z)|^{2}=d^{\prime} d^{\prime \prime} \log \left(1+|z|^{2}\right) \quad \text { on } \mathbb{C}^{n},  \tag{15.8}\\
\Theta(\mathscr{O}(1))_{a} & =\sum_{1 \leqslant j \leqslant n} d z_{j} \wedge d \bar{z}_{j} .
\end{align*}
$$

On the other hand, Th. 14.3 and (14.7) imply

$$
d^{\prime \prime} g^{\star}=-j \circ \beta^{\star}, \quad \Theta(H)=\beta \wedge \beta^{\star},
$$

where $j: \mathscr{O}(-1) \longrightarrow \underline{V}$ is the inclusion, $g^{\star}: H \longrightarrow \underline{V}$ the orthogonal splitting and $\beta^{\star} \in$ $\mathscr{C}_{0,1}^{\infty}(P(V), \operatorname{Hom}(H, \mathscr{O}(-1)))$. The images $\left(\widetilde{e}_{1}, \ldots, \widetilde{e}_{n}\right)$ of $e_{1}, \ldots, e_{n}$ in $H=\underline{V} / \mathbb{O}(-1)$ define a holomorphic frame of $H_{\mid \mathbb{C}^{n}}$ and we have

$$
\begin{align*}
g^{\star} \cdot \widetilde{e}_{j} & =e_{j}-\frac{\left\langle e_{j}, \varepsilon\right\rangle}{|\varepsilon|^{2}}=e_{j}-\frac{\bar{z}_{j}}{1+|z|^{2}} \varepsilon, \quad d^{\prime \prime} g_{a}^{\star} \cdot \widetilde{e}_{j}=-d \bar{z}_{j} \otimes \varepsilon \\
\beta_{a}^{\star} & =\sum_{1 \leqslant j \leqslant n} d \bar{z}_{j} \otimes \widetilde{e}_{j}^{\star} \otimes \varepsilon, \quad \beta_{a}=\sum_{1 \leqslant j \leqslant n} d z_{j} \otimes \varepsilon^{\star} \otimes \widetilde{e}_{j} \\
\Theta(H)_{a} & =\sum_{1 \leqslant j, k \leqslant n} d z_{j} \wedge d \bar{z}_{k} \otimes \widetilde{e}_{k}^{\star} \otimes \widetilde{e}_{j} \tag{15.9}
\end{align*}
$$

(15.10) Theorem. The cohomology algebra $H^{\bullet}\left(\mathbb{P}^{n}, \mathbb{Z}\right)$ is isomorphic to the quotient ring $\mathbb{Z}[h] /\left(h^{n+1}\right)$ where the generator $h$ is given by $h=c_{1}(\mathbb{G}(1))$ in $H^{2}\left(\mathbb{P}^{n}, \mathbb{Z}\right)$.

Proof. Consider the inclusion $\mathbb{P}^{n-1}=P\left(\mathbb{C}^{n} \times\{0\}\right) \subset \mathbb{P}^{n}$. Topologically, $\mathbb{P}^{n}$ is obtained from $\mathbb{P}^{n-1}$ by attaching a $2 n$-cell $B_{2 n}$ to $\mathbb{P}^{n-1}$, via the map

$$
\begin{aligned}
f: B_{2 n} & \longrightarrow \mathbb{P}^{n} \\
z & \longmapsto\left[z, 1-|z|^{2}\right], \quad z \in \mathbb{C}^{n}, \quad|z| \leqslant 1
\end{aligned}
$$

which sends $S^{2 n-1}=\{|z|=1\}$ onto $\mathbb{P}^{n-1}$. That is, $\mathbb{P}^{n}$ is homeomorphic to the quotient space of $B_{2 n} \amalg \mathbb{P}^{n-1}$, where every point $z \in S^{2 n-1}$ is identified with its image $f(z) \in \mathbb{P}^{n-1}$. We shall prove by induction on $n$ that

$$
\begin{equation*}
H^{2 k}\left(\mathbb{P}^{n}, \mathbb{Z}\right)=\mathbb{Z}, \quad 0 \leqslant k \leqslant n, \quad \text { otherwise } \quad H^{l}\left(\mathbb{P}^{n}, \mathbb{Z}\right)=0 \tag{15.11}
\end{equation*}
$$

The result is clear for $\mathbb{P}^{0}$, which is reduced to a single point. For $n \geqslant 1$, consider the covering $\left(U_{1}, U_{2}\right)$ of $\mathbb{P}^{n}$ such that $U_{1}$ is the image by $f$ of the open ball $B_{2 n}^{\circ}$ and $U_{2}=\mathbb{P}^{n} \backslash\{f(0)\}$. Then $U_{1} \approx B_{2 n}^{\circ}$ is contractible, whereas $U_{2}=\left(B_{2 n} \backslash\{0\}\right) \amalg_{S^{2 n-1}} \mathbb{P}^{n-1}$. Moreover $U_{1} \cap U_{2} \approx B_{2 n}^{\circ} \backslash\{0\}$ can be retracted on the $(2 n-1)$-sphere of radius $1 / 2$. For $q \geqslant 2$, the Mayer-Vietoris exact sequence IV-3.11 yields

$$
\begin{aligned}
& \cdots H^{q-1}\left(\mathbb{P}^{n-1}, \mathbb{Z}\right) \\
& \longrightarrow H^{q-1}\left(S^{2 n-1}, \mathbb{Z}\right) \\
& H^{q}\left(\mathbb{P}^{n}, \mathbb{Z}\right) \longrightarrow H^{q}\left(\mathbb{P}^{n-1}, \mathbb{Z}\right) \longrightarrow H^{q}\left(S^{2 n-1}, \mathbb{Z}\right) \cdots
\end{aligned}
$$

For $q=1$, the first term has to be replaced by $H^{0}\left(\mathbb{P}^{n-1}, \mathbb{Z}\right) \oplus \mathbb{Z}$, so that the first arrow is onto. Formula (15.11) follows easily by induction, thanks to our computation of the cohomology groups of spheres in IV-14.6.

We know that $h=c_{1}(\mathbb{O}(1)) \in H^{2}\left(\mathbb{P}^{n}, \mathbb{Z}\right)$. It will follow necessarily that $h^{k}$ is a generator of $H^{2 k}\left(\mathbb{P}^{n}, \mathbb{Z}\right)$ if we can prove that $h^{n}$ is the fundamental class in $H^{2 n}(\mathbb{P}, \mathbb{Z})$, or equivalently that

$$
\begin{equation*}
c_{1}(\mathscr{O}(1))_{\mathbb{R}^{n}}^{n}=\int_{\mathbb{P}^{n}}\left(\frac{\mathrm{i}}{2 \pi} \Theta(\mathscr{O}(1))\right)^{n}=1 . \tag{15.12}
\end{equation*}
$$

This equality can be verified directly by means of (15.8), but we will avoid this computation. Observe that the element $e_{n}^{\star} \in\left(\mathbb{C}^{n+1}\right)^{\star}$ defines a section $\widetilde{e}_{n}^{\star}$ of $H^{0}\left(\mathbb{P}^{n}, \mathbb{O}(1)\right)$ transverse to 0 , whose zero set is the hyperplane $\mathbb{P}^{n-1}$. As $\left\{\frac{i}{2 \pi} \Theta(\mathscr{O}(1))\right\}=\left\{\left[\mathbb{P}^{n-1}\right]\right\}$ by Th. 13.4, we get

$$
\begin{aligned}
c_{1}(\overparen{O}(1)) & =\int_{\mathbb{P}^{1}}\left[\mathbb{P}^{0}\right]=1 \quad \text { for } n=1 \text { and } \\
c_{1}(\mathscr{O}(1))^{n} & =\int_{\mathbb{P}^{n}}\left[\mathbb{P}^{n-1}\right] \wedge\left(\frac{\mathrm{i}}{2 \pi} \Theta(\mathscr{O}(1))\right)^{n-1}=\int_{\mathbb{P}^{n-1}}\left(\frac{\mathrm{i}}{2 \pi} \Theta(\mathscr{O}(1))\right)^{n-1}
\end{aligned}
$$

in general. Since $\mathfrak{G}(-1)_{\mid \mathbb{P}^{n-1}}$ can be identified with the tautological line subbundle $\mathscr{O}_{\mathbb{P}^{n-1}}(-1)$ over $\mathbb{P}^{n-1}$, we have $\Theta(\mathscr{G}(1))_{\mathbb{P}^{n-1}}=\Theta\left(\mathbb{O}_{\mathbb{P}^{n-1}}(1)\right)$ and the proof is achieved by induction on $n$.

## § 15.C. Tautological Line Bundle Associated to a Vector Bundle

Let $E$ be a holomorphic vector bundle of rank $r$ over a complex manifold $X$. The projectivized bundle $P(E)$ is the bundle with $\mathbb{P}^{r-1}$ fibers over $X$ defined by $P(E)_{x}=P\left(E_{x}\right)$ for all $x \in X$. The points of $P(E)$ can thus be identified with the lines in the fibers of $E$. For any trivialization $\theta_{\alpha}: E_{\mid U_{\alpha}} \rightarrow U_{\alpha} \times \mathbb{C}^{r}$ of $E$ we have a corresponding trivialization $\widetilde{\theta}_{\alpha}: P(E)_{\mid U_{\alpha}} \rightarrow U_{\alpha} \times \mathbb{P}^{r-1}$, and it is clear that the transition automorphisms are the projectivizations $\widetilde{g}_{\alpha \beta} \in H^{0}\left(U_{\alpha} \cap U_{\beta}, P G L(r, \mathbb{C})\right)$ of the transition automorphisms $g_{\alpha \beta}$ of $E$.

Similarly, we have a dual projectivized bundle $P\left(E^{\star}\right)$ whose points can be identified with the hyperplanes of $E$ (every hyperplane $F$ in $E_{x}$ corresponds bijectively to the line of linear forms in $E_{x}^{\star}$ which vanish on $\left.F\right)$; note that $P(E)$ and $P\left(E^{\star}\right)$ coincide only when $r=\mathrm{rk} E=2$. If $\pi: P\left(E^{\star}\right) \rightarrow X$ is the natural projection, there is a tautological hyperplane subbundle $S$ of $\pi^{\star} E$ over $P\left(E^{\star}\right)$ such that $S_{[\xi]}=\xi^{-1}(0) \subset E_{x}$ for all $\xi \in E_{x}^{\star} \backslash\{0\}$.
[exercise: check that $S$ is actually locally trivial over $P\left(E^{\star}\right)$ ].
(15.13) Definition. The quotient line bundle $\pi^{\star} E / S$ is denoted $\mathscr{O}_{E}(1)$ and is called the tautological line bundle associated to $E$. Hence there is an exact sequence

$$
0 \longrightarrow S \longrightarrow \pi^{\star} E \longrightarrow \mathscr{O}_{E}(1) \longrightarrow 0
$$

of vector bundles over $P\left(E^{\star}\right)$.

Note that (13.3) applied with $V=E_{x}^{\star}$ implies that the restriction of $\mathscr{O}_{E}(1)$ to each fiber $P\left(E_{x}^{\star}\right) \simeq \mathbb{P}^{r-1}$ coincides with the line bundle $\mathscr{O}(1)$ introduced in Def. 15.2. Theorem 15.5 can then be extended to the present situation and yields:
(15.14) Theorem. For every $k \in \mathbb{Z}$, the direct image sheaf $\pi_{\star} \mathscr{O}_{E}(k)$ on $X$ vanishes for $k<0$ and is isomorphic to $\mathfrak{G}\left(S^{k} E\right)$ for $k \geqslant 0$.

Proof. For $k \geqslant 0$, the $k$-th symmetric power of the morphism $\pi^{\star} E \rightarrow \mathscr{O}_{E}(1)$ gives a morphism $\pi^{\star} S^{k} E \rightarrow \mathscr{O}_{E}(k)$. This morphism together with the pull-back morphism yield canonical arrows

$$
\Phi_{U}: H^{0}\left(U, S^{k} E\right) \xrightarrow{\pi^{\star}} H^{0}\left(\pi^{-1}(U), \pi^{\star} S^{k} E\right) \longrightarrow H^{0}\left(\pi^{-1}(U), \mathscr{O}_{E}(k)\right)
$$

for any open set $U \subset X$. The right hand side is by definition the space of sections of $\pi_{\star} \mathscr{O}_{E}(k)$ over $U$, hence we get a canonical sheaf morphism

$$
\Phi: \mathscr{O}\left(S^{k} E\right) \longrightarrow \pi_{\star} \mathscr{O}_{E}(k) .
$$

It is easy to check that this $\Phi$ coincides with the map $\Phi$ introduced in the proof of Cor. 15.6 when $X$ is reduced to a point. In order to check that $\Phi$ is an isomorphism, we may suppose that $U$ is chosen so small that $E_{\uparrow U}$ is trivial, say $E_{\lceil U}=U \times V$ with $\operatorname{dim} V=r$. Then $P\left(E^{\star}\right)=U \times P\left(V^{\star}\right)$ and $\mathscr{O}_{E}(1)=p^{\star} \mathscr{O}(1)$ where $\mathscr{O}(1)$ is the tautological line bundle over $P\left(V^{\star}\right)$ and $p: P\left(E^{\star}\right) \rightarrow P\left(V^{\star}\right)$ is the second projection. Hence we get

$$
\begin{aligned}
H^{0}\left(\pi^{-1}(U), \mathscr{O}_{E}(k)\right) & =H^{0}\left(U \times P\left(V^{\star}\right), p^{\star} \mathfrak{O}(1)\right) \\
& =\mathscr{O}_{X}(U) \otimes H^{0}\left(P\left(V^{\star}\right), \mathscr{O}(1)\right) \\
& =\mathscr{O}_{X}(U) \otimes S^{k} V=H^{0}\left(U, S^{k} E\right),
\end{aligned}
$$

as desired; the reason for the second equality is that $p^{\star} \mathscr{O}(1)$ coincides with $\mathfrak{G}(1)$ on each fiber $\{x\} \times P\left(V^{\star}\right)$ of $p$, thus any section of $p^{\star} \mathscr{O}(1)$ over $U \times P\left(V^{\star}\right)$ yields a family of sections $H^{0}\left(\{x\} \times P\left(V^{\star}\right), \mathscr{O}(k)\right)$ depending holomorphically in $x$. When $k<0$ there are no non zero such sections, thus $\pi_{\star} \mathfrak{O}_{E}(k)=0$.

Finally, suppose that $E$ is equipped with a hermitian metric. Then the morphism $\pi^{\star} E \rightarrow \mathfrak{O}_{E}(1)$ endows $\mathscr{O}_{E}(1)$ with a quotient metric. We are going to compute the associated curvature form $\Theta\left(\mathbb{Q}_{E}(1)\right)$.

Fix a point $x_{0} \in X$ and $a \in P\left(E_{x_{0}}^{\star}\right)$. Then Prop. 12.10 implies the existence of a normal coordinate frame $\left.\left(e_{\lambda}\right)_{1 \leqslant \lambda \leqslant r}\right)$ of $E$ at $x_{0}$ such that $a$ is the hyperplane $\left\langle e_{2}, \ldots, e_{r}\right\rangle=\left(e_{1}^{\star}\right)^{-1}(0)$ at $x_{0}$. Let $\left(z_{1}, \ldots, z_{n}\right)$ be local coordinates on $X$ near $x_{0}$ and let $\left(\xi_{1}, \ldots, \xi_{r}\right)$ be coordinates on $E^{\star}$ with respect to the dual frame $\left(e_{1}^{\star}, \ldots, e_{r}^{\star}\right)$. If we assign $\xi_{1}=1$, then $\left(z_{1}, \ldots, z_{n}, \xi_{2}, \ldots, \xi_{r}\right)$ define local coordinates on $P\left(E^{\star}\right)$ near $a$, and we have a local section of $\mathscr{O}_{E}(-1):=\mathscr{O}_{E}(1)^{\star} \subset \pi^{\star} E^{\star}$ defined by

$$
\varepsilon(z, \xi)=e_{1}^{\star}(z)+\sum_{2 \leqslant \lambda \leqslant r} \xi_{\lambda} e_{\lambda}^{\star}(z) .
$$

The hermitian matrix $\left(\left\langle e_{\lambda}^{\star}, e_{\mu}^{\star}\right\rangle\right)$ is just the congugate of the inverse of matrix $\left(\left\langle e_{\lambda}, e_{\mu}\right\rangle\right)=$ $\operatorname{Id}-\left(\sum c_{j k \lambda \mu} z_{j} \bar{z}_{k}\right)+O\left(|z|^{3}\right)$, hence we get

$$
\left\langle e_{\lambda}^{\star}(z), e_{\mu}^{\star}(z)\right\rangle=\delta_{\lambda \mu}+\sum_{1 \leqslant j, k \leqslant n} c_{j k \mu \lambda} z_{j} \bar{z}_{k}+O\left(|z|^{3}\right),
$$

where ( $c_{j k \lambda \mu}$ ) are the curvature coefficients of $\Theta(E)$; accordingly we have $\Theta\left(E^{\star}\right)=$ $-\Theta(E)^{\dagger}$. We infer from this

$$
|\varepsilon(z, \xi)|^{2}=1+\sum_{1 \leqslant j, k \leqslant n} c_{j k 11} z_{j} \bar{z}_{k}+\sum_{2 \leqslant \lambda \leqslant r}\left|\xi_{\lambda}\right|^{2}+O\left(|z|^{3}\right) .
$$

Since $\Theta\left(\mathscr{O}_{E}(1)\right)=d^{\prime} d^{\prime \prime} \log |\varepsilon(z, \xi)|^{2}$, we get

$$
\Theta\left(\mathscr{O}_{E}(1)\right)_{a}=\sum_{1 \leqslant j, k \leqslant n} c_{j k 11} d z_{j} \wedge d \bar{z}_{k}+\sum_{2 \leqslant \lambda \leqslant r} d \xi_{\lambda} \wedge d \bar{\xi}_{\lambda} .
$$

Note that the first summation is simply $-\left\langle\Theta\left(E^{\star}\right) a, a\right\rangle /|a|^{2}=-$ curvature of $E^{\star}$ in the direction $a$. A unitary change of variables then gives the slightly more general formula:
(15.15) Formula. Let $\left(e_{\lambda}\right)$ be a normal coordinate frame of $E$ at $x_{0} \in X$ and let $\Theta(E)_{x_{0}}=\sum c_{j k \lambda \mu} d z_{j} \wedge d \bar{z}_{k} \otimes e_{\lambda}^{\star} \otimes e_{\mu}$. At any point $a \in P\left(E^{\star}\right)$ represented by a vector $\sum a_{\lambda} e_{\lambda}^{\star} \in E_{x_{0}}^{\star}$ of norm 1 , the curvature of $\mathscr{O}_{E}(1)$ is

$$
\Theta\left(\mathbb{O}_{E}(1)\right)_{a}=\sum_{1 \leqslant j, k \leqslant n, 1 \leqslant \lambda, \mu \leqslant r} c_{j k \mu \lambda} a_{\lambda} \bar{a}_{\mu} d z_{j} \wedge d \bar{z}_{k}+\sum_{1 \leqslant \lambda \leqslant r-1} d \zeta_{\lambda} \wedge d \bar{\zeta}_{\lambda},
$$

where $\left(\zeta_{\lambda}\right)$ are coordinates near a on $P\left(E^{\star}\right)$, induced by unitary coordinates on the hyperplane $a^{\perp} \subset E_{x_{0}}^{\star}$.

## § 16. Grassmannians and Universal Vector Bundles

## $\S$ 16.A. Universal Subbundles and Quotient Vector Bundles

If $V$ is a complex vector space of dimension $d$, we denote by $G_{r}(V)$ the set of all $r$-codimensional vector subspaces of $V$. Let $a \in G_{r}(V)$ and $W \subset V$ be fixed such that

$$
V=a \oplus W, \quad \operatorname{dim}_{\mathbb{C}} W=r
$$

Then any subspace $x \in G_{r}(V)$ in the open subset

$$
\Omega_{W}=\left\{x \in G_{r}(V) ; x \oplus W=V\right\}
$$

can be represented in a unique way as the graph of a linear map $u$ in $\operatorname{Hom}(a, W)$. This gives rise to a covering of $G_{r}(V)$ by affine coordinate charts $\Omega_{W} \simeq \operatorname{Hom}(a, W) \simeq \mathbb{C}^{r(d-r)}$. Indeed, let $\left(e_{1}, \ldots, e_{r}\right)$ and $\left(e_{r+1}, \ldots, e_{n}\right)$ be respective bases of $W$ and $a$. Every point $x \in \Omega_{W}$ is the graph of a linear map

$$
\begin{equation*}
u: a \longrightarrow W, \quad u\left(e_{k}\right)=\sum_{1 \leqslant j \leqslant r} z_{j k} e_{j}, \quad r+1 \leqslant k \leqslant d, \tag{16.1}
\end{equation*}
$$

i.e. $x=\operatorname{Vect}\left(e_{k}+\sum_{1 \leqslant j \leqslant r} z_{j k} e_{j}\right)_{r+1 \leqslant k \leqslant d}$. We choose $\left(z_{j k}\right)$ as complex coordinates on $\Omega_{W}$. These coordinates are centered at $a=\operatorname{Vect}\left(e_{r+1}, \ldots, e_{d}\right)$.
(16.2) Proposition. $G_{r}(V)$ is a compact complex analytic manifold of dimension $n=$ $r(d-r)$.

Proof. It is immediate to verify that the coordinate change between two affine charts of $G_{r}(V)$ is holomorphic. Fix an arbitrary hermitian metric on $V$. Then the unitary group $U(V)$ is compact and acts transitively on $G_{r}(V)$. The isotropy subgroup of a point $a \in G_{r}(V)$ is $U(a) \times U\left(a^{\perp}\right)$, hence $G_{r}(V)$ is diffeomorphic to the compact quotient space $U(V) / U(a) \times U\left(a^{\perp}\right)$.

Next, we consider the tautological subbundle $S \subset \underline{V}:=G_{r}(V) \times V$ defined by $S_{x}=x$ for all $x \in G_{r}(V)$, and the quotient bundle $Q=\underline{V} / S$ of rank $r$ :

$$
\begin{equation*}
0 \longrightarrow S \longrightarrow \underline{V} \longrightarrow Q \longrightarrow 0 \tag{16.3}
\end{equation*}
$$

An interesting special case is $r=d-1, G_{d-1}(V)=P(V), S=\mathscr{O}(-1), Q=H$. The case $r=1$ is dual, we have the identification $G_{1}(V)=P\left(V^{\star}\right)$ because every hyperplane $x \subset V$ corresponds bijectively to the line in $V^{\star}$ of linear forms $\xi \in V^{\star}$ that vanish on $x$. Then the bundles $\mathfrak{G}(-1) \subset \underline{V}^{\star}$ and $H$ on $P\left(V^{\star}\right)$ are given by

$$
\begin{aligned}
\mathfrak{O}(-1)_{[\xi]} & =\mathbb{C} \cdot \xi \simeq(V / x)^{\star}=Q_{x}^{\star}, \\
H_{[\xi]} & =V^{\star} / \mathbb{C} \cdot \xi \simeq x^{\star}=S_{x}^{\star},
\end{aligned}
$$

therefore $S=H^{\star}, Q=\mathscr{O}(1)$. This special case will allow us to compute $H^{0}\left(G_{r}(V), Q\right)$ in general.
(16.4) Proposition. There is an isomorphism

$$
V=H^{0}\left(G_{r}(V), \underline{V}\right) \xrightarrow{\sim} H^{0}\left(G_{r}(V), Q\right) .
$$

Proof. Let $V=W \oplus W^{\prime}$ be an arbitrary direct sum decomposition of $V$ with $\operatorname{codim} W=$ $r-1$. Consider the projective space

$$
P\left(W^{\star}\right)=G_{1}(W) \subset G_{r}(V),
$$

its tautological hyperplane subbundle $H^{\star} \subset \underline{W}=P\left(W^{\star}\right) \times W$ and the exact sequence $0 \rightarrow H^{\star} \rightarrow \underline{W} \rightarrow \mathbb{O}(1) \rightarrow 0$. Then $S_{\upharpoonright P\left(W^{\star}\right)}$ coincides with $H^{\star}$ and

$$
Q_{\upharpoonright P\left(W^{\star}\right)}=\left(\underline{W} \oplus \underline{W}^{\prime}\right) / H^{\star}=\left(\underline{W} / H^{\star}\right) \oplus \underline{W}^{\prime}=\overparen{O}(1) \oplus \underline{W}^{\prime} .
$$

Theorem 15.5 implies $H^{0}\left(P\left(W^{\star}\right), \mathscr{O}(1)\right)=W$, therefore the space

$$
H^{0}\left(P\left(W^{\star}\right), Q_{\mid P\left(W^{\star}\right)}\right)=W \oplus W^{\prime}
$$

is generated by the images of the constant sections of $\underline{V}$. Since $W$ is arbitrary, Prop. 16.4 follows immediately.

Let us compute the tangent space $T G_{r}(V)$. The linear group $\mathrm{Gl}(V)$ acts transitively on $G_{r}(V)$, and the tangent space to the isotropy subgroup of a point $x \in G_{r}(V)$ is the set of elements $u \in \operatorname{Hom}(V, V)$ in the Lie algebra such that $u(x) \subset x$. We get therefore

$$
\begin{aligned}
T_{x} G_{r}(V) & \simeq \operatorname{Hom}(V, V) /\{u ; u(x) \subset x\} \\
& \simeq \operatorname{Hom}(V, V / x) /\{\widetilde{u} ; \widetilde{u}(x)=\{0\}\} \\
& \simeq \operatorname{Hom}(x, V / x)=\operatorname{Hom}\left(S_{x}, Q_{x}\right)
\end{aligned}
$$

(16.5) Corollary. $T G_{r}(V)=\operatorname{Hom}(S, Q)=S^{\star} \otimes Q$.

## § 16.B. Plücker Embedding

There is a natural map, called the Plücker embedding,

$$
\begin{equation*}
j_{r}: G_{r}(V) \hookrightarrow P\left(\Lambda^{r} V^{\star}\right) \tag{16.6}
\end{equation*}
$$

constructed as follows. If $x \in G_{r}(V)$ is defined by $r$ independent linear forms $\xi_{1}, \ldots, \xi_{r} \in$ $V^{\star}$, we set

$$
j_{r}(x)=\left[\xi_{1} \wedge \cdots \wedge \xi_{r}\right] .
$$

Then $x$ is the subspace of vectors $v \in V$ such that $v\lrcorner\left(\xi_{1} \wedge \cdots \wedge \xi_{r}\right)=0$, so $j_{r}$ is injective. Since the linear group $\mathrm{Gl}(V)$ acts transitively on $G_{r}(V)$, the rank of the differential $d j_{r}$ is a constant. As $j_{r}$ is injective, the constant rank theorem implies:
(16.7) Proposition. The map $j_{r}$ is a holomorphic embedding.

Now, we define a commutative diagram

as follows: for $x=\xi_{1}^{-1}(0) \cap \cdots \cap \xi_{r}^{-1}(0) \in G_{r}(V)$ and $\widetilde{v}=\widetilde{v}_{1} \wedge \cdots \wedge \widetilde{v}_{r} \in \Lambda^{r} Q_{x}$ where $\widetilde{v}_{k} \in Q_{x}=V / x$ is the image of $v_{k} \in V$ in the quotient, we let $J_{r}(\widetilde{v}) \in \mathscr{O}(1)_{j_{r}(x)}$ be the linear form on $\mathfrak{G}(-1)_{j_{r}(x)}=\mathbb{C} . \xi_{1} \wedge \ldots \wedge \xi_{r}$ such that

$$
\left\langle J_{r}(\widetilde{v}), \lambda \xi_{1} \wedge \ldots \wedge \xi_{r}\right\rangle=\lambda \operatorname{det}\left(\xi_{j}\left(v_{k}\right)\right), \quad \lambda \in \mathbb{C}
$$

Then $J_{r}$ is an isomorphism on the fibers, so $\Lambda^{r} Q$ can be identified with the pull-back of ©(1) by $j_{r}$.

## § 16.C. Curvature of the Universal Vector Bundles

Assume now that $V$ is a hermitian vector space. We shall generalize our curvature computations of $\S 15 . \mathrm{C}$ to the present situation. Let $a \in G_{r}(V)$ be a given point. We take $W$ to be the orthogonal complement of $a$ in $V$ and select an orthonormal basis $\left(e_{1}, \ldots, e_{d}\right)$ of $V$ such that $W=\operatorname{Vect}\left(e_{1}, \ldots, e_{r}\right), a=\operatorname{Vect}\left(e_{r+1}, \ldots, e_{d}\right)$. For any point $x \in G_{r}(V)$ in $\Omega_{W}$ with coordinates $\left(z_{j k}\right)$, we set

$$
\begin{aligned}
& \varepsilon_{k}(x)=e_{k}+\sum_{1 \leqslant j \leqslant r} z_{j k} e_{j}, \quad r+1 \leqslant k \leqslant d, \\
& \widetilde{e}_{j}(x)=\text { image of } e_{j} \text { in } Q_{x}=V / x, \quad 1 \leqslant j \leqslant r .
\end{aligned}
$$

Then $\left(\widetilde{e}_{1}, \ldots, \widetilde{e}_{r}\right)$ and $\left(\varepsilon_{r+1}, \ldots, \varepsilon_{d}\right)$ are holomorphic frames of $Q$ and $S$ respectively. If $g^{\star}: Q \longrightarrow \underline{V}$ is the orthogonal splitting of $g: \underline{V} \longrightarrow Q$, then

$$
g^{\star} \cdot \widetilde{e}_{j}=e_{j}+\sum_{r+1 \leqslant k \leqslant d} \zeta_{j k} \varepsilon_{k}
$$

for some $\zeta_{j k} \in \mathbb{C}$. After an easy computation we find

$$
0=\left\langle\widetilde{e}_{j}, g \varepsilon_{k}\right\rangle=\left\langle g^{\star} \widetilde{e}_{j}, \varepsilon_{k}\right\rangle=\zeta_{j k}+\bar{z}_{j k}+\sum_{l, m} \zeta_{j m} z_{l m} \bar{z}_{l k},
$$

so that $\zeta_{j k}=-\bar{z}_{j k}+O\left(|z|^{2}\right)$. Formula (13.3) yields

$$
\begin{align*}
d^{\prime \prime} g_{a}^{\star} \cdot \widetilde{e}_{j} & =-\sum_{r+1 \leqslant k \leqslant d} d \bar{z}_{j k} \otimes \varepsilon_{k}, \\
\beta_{a}^{\star} & =\sum_{j, k} d \bar{z}_{j k} \otimes \widetilde{e}_{j}^{\star} \otimes \varepsilon_{k}, \quad \beta_{a}=\sum_{j, k} d z_{j k} \otimes \varepsilon_{k}^{\star} \otimes \widetilde{e}_{j}, \\
\Theta(Q)_{a} & =\left(\beta \wedge \beta^{\star}\right)_{a}=\sum_{j, k, l} d z_{j k} \wedge d \bar{z}_{l k} \otimes \widetilde{e}_{l}^{\star} \otimes \widetilde{e}_{j},  \tag{16.9}\\
\Theta(S)_{a} & =\left(\beta^{\star} \wedge \beta\right)_{a}=-\sum_{j, k, l} d z_{j k} \wedge d \bar{z}_{j l} \otimes \varepsilon_{k}^{\star} \otimes \varepsilon_{l} . \tag{16.10}
\end{align*}
$$

## Chapter VI

## Hodge Theory

The goal of this chapter is to prove a number of basic facts in the Hodge theory of real or complex manifolds. The theory rests essentially on the fact that the De Rham (or Dolbeault) cohomology groups of a compact manifold can be represented by means of spaces of harmonic forms, once a Riemannian metric has been chosen. At this point, some knowledge of basic results about elliptic differential operators is required. The special properties of compact Kähler manifolds are then investigated in detail: Hodge decomposition theorem, hard Lefschetz theorem, Jacobian and Albanese variety, ...; the example of curves is treated in detail. Finally, the Hodge-Frölicher spectral sequence is applied to get some results on general compact complex manifolds, and it is shown that Hodge decomposition still holds for manifolds in the Fujiki class ( ${ }^{\mathscr{C}}$ ).

## § 1. Differential Operators on Vector Bundles

We first describe some basic concepts concerning differential operators (symbol, composition, adjunction, ellipticity), in the general setting of vector bundles. Let $M$ be a $\mathscr{C}^{\infty}$ differentiable manifold, $\operatorname{dim}_{\mathbb{R}} M=m$, and let $E, F$ be $\mathbb{K}$-vector bundles over $M$, with $\mathbb{K}=\mathbb{R}$ or $\mathbb{K}=\mathbb{C}, \operatorname{rank} E=r, \operatorname{rank} F=r^{\prime}$.
(1.1) Definition. A (linear) differential operator of degree $\delta$ from $E$ to $F$ is a $\mathbb{K}$-linear operator $P: \mathscr{C}^{\infty}(M, E) \rightarrow \mathscr{C}^{\infty}(M, F), u \mapsto P u$ of the form

$$
P u(x)=\sum_{|\alpha| \leqslant \delta} a_{\alpha}(x) D^{\alpha} u(x),
$$

where $E_{\uparrow \Omega} \simeq \Omega \times \mathbb{K}^{r}, F_{\upharpoonright \Omega} \simeq \Omega \times \mathbb{K}^{r^{\prime}}$ are trivialized locally on some open chart $\Omega \subset M$ equipped with local coordinates $\left(x_{1}, \ldots, x_{m}\right)$, and where $a_{\alpha}(x)=\left(a_{\alpha \lambda \mu}(x)\right)_{1 \leqslant \lambda \leqslant r^{\prime}, 1 \leqslant \mu \leqslant r}$ are $r^{\prime} \times r$-matrices with $\mathscr{C}^{\infty}$ coefficients on $\Omega$. Here $D^{\alpha}=\left(\partial / \partial x_{1}\right)^{\alpha_{1}} \ldots\left(\partial / \partial x_{m}\right)^{\alpha_{m}}$ as usual, and $u=\left(u_{\mu}\right)_{1 \leqslant \mu \leqslant r}, D^{\alpha} u=\left(D^{\alpha} u_{\mu}\right)_{1 \leqslant \mu \leqslant r}$ are viewed as column matrices.

If $t \in \mathbb{K}$ is a parameter, a simple calculation shows that $e^{-t u(x)} P\left(e^{t u(x)}\right)$ is a polynomial of degree $\delta$ in $t$, of the form

$$
e^{-t u(x)} P\left(e^{t u(x)}\right)=t^{\delta} \sigma_{P}(x, d u(x))+\text { lower order terms } c_{j}(x) t^{j}, j<\delta,
$$

where $\sigma_{P}$ is the polynomial map from $T_{M}^{\star} \rightarrow \operatorname{Hom}(E, F)$ defined by

$$
\begin{equation*}
T_{M, x}^{\star} \ni \xi \mapsto \sigma_{P}(x, \xi) \in \operatorname{Hom}\left(E_{x}, F_{x}\right), \quad \sigma_{P}(x, \xi)=\sum_{|\alpha|=\delta} a_{\alpha}(x) \xi^{\alpha} \tag{1.2}
\end{equation*}
$$

The formula involving $e^{-t u} P\left(e^{t u}\right)$ shows that $\sigma_{P}(x, \xi)$ actually does not depend on the choice of coordinates nor on the trivializations used for $E, F$. It is clear that $\sigma_{P}(x, \xi)$ is smooth on $T_{M}^{\star}$ as a function of $(x, \xi)$, and is a homogeneous polynomial of degree $\delta$ in $\xi$. We say that $\sigma_{P}$ is the principal symbol of $P$. Now, if $E, F, G$ are vector bundles and

$$
P: \mathscr{C}^{\infty}(M, E) \rightarrow \mathscr{C}^{\infty}(M, F), \quad Q: \mathscr{C}^{\infty}(M, F) \rightarrow \mathscr{C}^{\infty}(M, G)
$$

are differential operators of respective degrees $\delta_{P}, \delta_{Q}$, it is easy to check that $Q \circ P$ : $\mathscr{C}^{\infty}(M, E) \rightarrow \mathscr{C}^{\infty}(M, G)$ is a differential operator of degree $\delta_{P}+\delta_{Q}$ and that

$$
\begin{equation*}
\sigma_{Q \circ P}(x, \xi)=\sigma_{Q}(x, \xi) \sigma_{P}(x, \xi) \tag{1.3}
\end{equation*}
$$

Here the product of symbols is computed as a product of matrices.
Now, assume that $M$ is oriented and is equipped with a smooth volume form $d V(x)=$ $\gamma(x) d x_{1} \wedge \ldots d x_{m}$, where $\gamma(x)>0$ is a smooth density. If $E$ is a euclidean or hermitian vector bundle, we have a Hilbert space $L^{2}(M, E)$ of global sections $u$ of $E$ with measurable coefficients, satisfying the $L^{2}$ estimate

$$
\begin{equation*}
\|u\|^{2}=\int_{M}|u(x)|^{2} d V(x)<+\infty . \tag{1.4}
\end{equation*}
$$

We denote by

$$
\langle\langle u, v\rangle\rangle=\int_{M}\langle u(x), v(x)\rangle d V(x), \quad u, v \in L^{2}(M, E)
$$

the corresponding $L^{2}$ inner product.
(1.5) Definition. If $P: \mathscr{C}^{\infty}(M, E) \rightarrow \mathscr{C}^{\infty}(M, F)$ is a differential operator and both $E$, $F$ are euclidean or hermitian, there exists a unique differential operator

$$
P^{\star}: \mathscr{C}^{\infty}(M, F) \rightarrow \mathscr{C}^{\infty}(M, E),
$$

called the formal adjoint of $P$, such that for all sections $u \in \mathscr{C}^{\infty}(M, E)$ and $v \in$ $C^{\infty}(M, F)$ there is an identity

$$
\langle\langle P u, v\rangle\rangle=\left\langle\left\langle u, P^{\star} v\right\rangle\right\rangle, \quad \text { whenever Supp } u \cap \operatorname{Supp} v \subset \subset M .
$$

Proof. The uniqueness is easy, using the density of the set of elements $u \in \mathscr{C}^{\infty}(M, E)$ with compact support in $L^{2}(M, E)$. Since uniqueness is clear, it is enough, by a partition of unity argument, to show the existence of $P^{\star}$ locally. Now, let $P u(x)=$ $\sum_{|\alpha| \leqslant \delta} a_{\alpha}(x) D^{\alpha} u(x)$ be the expansion of $P$ with respect to trivializations of $E, F$ given by orthonormal frames over some coordinate open set $\Omega \subset M$. Assuming $\operatorname{Supp} u \cap$ Supp $v \subset \subset \Omega$, an integration by parts yields

$$
\begin{aligned}
\langle\langle P u, v\rangle\rangle & =\int_{\Omega} \sum_{|\alpha| \leqslant \delta, \lambda, \mu} a_{\alpha \lambda \mu} D^{\alpha} u_{\mu}(x) \bar{v}_{\lambda}(x) \gamma(x) d x_{1}, \ldots, d x_{m} \\
& =\int_{\Omega} \sum_{|\alpha| \leqslant \delta, \lambda, \mu}(-1)^{|\alpha|} u_{\mu}(x) \overline{D^{\alpha}\left(\gamma(x) \bar{a}_{\alpha \lambda \mu} v_{\lambda}(x)\right.} d x_{1}, \ldots, d x_{m} \\
& =\int_{\Omega}\left\langle u, \sum_{|\alpha| \leqslant \delta}(-1)^{|\alpha|} \gamma(x)^{-1} D^{\alpha}\left(\gamma(x)^{t} \bar{a}_{\alpha} v(x)\right)\right\rangle d V(x) .
\end{aligned}
$$

Hence we see that $P^{\star}$ exists and is uniquely defined by

$$
\begin{equation*}
P^{\star} v(x)=\sum_{|\alpha| \leqslant \delta}(-1)^{|\alpha|} \gamma(x)^{-1} D^{\alpha}\left(\gamma(x)^{t} \bar{a}_{\alpha} v(x)\right) . \tag{1.6}
\end{equation*}
$$

It follows immediately from (1.6) that the principal symbol of $P^{\star}$ is

$$
\begin{equation*}
\sigma_{P^{\star}}(x, \xi)=(-1)^{\delta} \sum_{|\alpha|=\delta}{ }^{t} \bar{a}_{\alpha} \xi^{\alpha}=(-1)^{\delta} \sigma_{P}(x, \xi)^{\star} \tag{1.7}
\end{equation*}
$$

(1.8) Definition. A differential operator $P$ is said to be elliptic if

$$
\sigma_{P}(x, \xi) \in \operatorname{Hom}\left(E_{x}, F_{x}\right)
$$

is injective for every $x \in M$ and $\xi \in T_{M, x}^{\star} \backslash\{0\}$.

## § 2. Formalism of PseudoDifferential Operators

We assume throughout this section that $(M, g)$ is a compact Riemannian manifold. For any positive integer $k$ and any hermitian bundle $F \rightarrow M$, we denote by $W^{k}(M, F)$ the Sobolev space of sections $s: M \rightarrow F$ whose derivatives up to order $k$ are in $L^{2}$. Let $\left\|\|_{k}\right.$ be the norm of the Hilbert space $W^{k}(M, F)$. Let $P$ be an elliptic differential operator of order $d$ acting on $\mathscr{C}^{\infty}(M, F)$. We need the following basic facts of elliptic $P D E$ theory, see e.g. [Hörmander 1963].
(2.1) Sobolev lemma. For $k>l+\frac{m}{2}, W^{k}(M, F) \subset C^{l}(M, F)$.
(2.2) Rellich lemma. For every integer $k$, the inclusion

$$
W^{k+1}(M, F) \hookrightarrow W^{k}(M, F)
$$

is a compact linear operator.
(2.3) Gårding's inequality. Let $\widetilde{P}$ be the extension of $P$ to sections with distribution coefficients. For any $u \in W^{0}(M, F)$ such that $\widetilde{P} u \in W^{k}(M, F)$, then $u \in W^{k+d}(M, F)$ and

$$
\|u\|_{k+d} \leqslant C_{k}\left(\|\widetilde{P} u\|_{k}+\|u\|_{0}\right),
$$

where $C_{k}$ is a positive constant depending only on $k$.
(2.4) Corollary. The operator $P: \mathscr{C}^{\infty}(M, F) \rightarrow \mathscr{C}^{\infty}(M, F)$ has the following properties:
i) $\operatorname{ker} P$ is finite dimensional.
ii) $P\left(\mathscr{C}^{\infty}(M, F)\right)$ is closed and of finite codimension; furthermore, if $P^{\star}$ is the formal adjoint of $P$, there is a decomposition

$$
\mathscr{C}^{\infty}(M, F)=P\left(\mathscr{C}^{\infty}(M, F)\right) \oplus \operatorname{ker} P^{\star}
$$

as an orthogonal direct sum in $W^{0}(M, F)=L^{2}(M, F)$.
Proof. (i) Gårding's inequality shows that $\|u\|_{k+d} \leqslant C_{k}\|u\|_{0}$ for any $u$ in ker $P$. Thanks to the Sobolev lemma, this implies that ker $P$ is closed in $W^{0}(M, F)$. Moreover, the unit closed $\left\|\|_{0}\right.$-ball of ker $P$ is contained in the $\| \|_{d}$-ball of radius $C_{0}$, thus compact by the Rellich lemma. Riesz' theorem implies that dim ker $P<+\infty$.
(ii) We first show that the extension

$$
\widetilde{P}: W^{k+d}(M, F) \rightarrow W^{k}(M, F)
$$

has a closed range for any $k$. For every $\varepsilon>0$, there exists a finite number of elements $v_{1}, \ldots, v_{N} \in W^{k+d}(M, F), N=N(\varepsilon)$, such that

$$
\begin{equation*}
\|u\|_{0} \leqslant \varepsilon\|u\|_{k+d}+\sum_{j=1}^{N}\left|\left\langle u, v_{j}\right\rangle_{0}\right| \tag{2.5}
\end{equation*}
$$

indeed the set

$$
K_{\left(v_{j}\right)}=\left\{u \in W^{k+d}(M, F) ; \varepsilon\|u\|_{k+d}+\sum_{j=1}^{N}\left|\left\langle\left\langle u, v_{j}\right\rangle\right\rangle_{0}\right| \leqslant 1\right\}
$$

is relatively compact in $W^{0}(M, F)$ and $\bigcap_{\left(v_{j}\right)} \bar{K}_{\left(v_{j}\right)}=\{0\}$. It follows that there exist elements $\left(v_{j}\right)$ such that $\bar{K}_{\left(v_{j}\right)}$ is contained in the unit ball of $W^{0}(M, F), Q E D$. Substitute $\|u\|_{0}$ by the upper bound (2.5) in Gårding's inequality; we get

$$
\left(1-C_{k} \varepsilon\right)\|u\|_{k+d} \leqslant C_{k}\left(\|\widetilde{P} u\|_{k}+\sum_{j=1}^{N} \mid\left\langle\left\langle u, v_{j}\right\rangle_{0}\right|\right)
$$

Define $G=\left\{u \in W^{k+d}(M, F) ; u \perp v_{j}, 1 \leqslant j \leqslant n\right\}$ and choose $\varepsilon=1 / 2 C_{k}$. We obtain

$$
\|u\|_{k+d} \leqslant 2 C_{k}\|\widetilde{P} u\|_{k}, \quad \forall u \in G
$$

This implies that $\widetilde{P}(G)$ is closed. Therefore

$$
\widetilde{P}\left(W^{k+d}(M, F)\right)=\widetilde{P}(G)+\operatorname{Vect}\left(\widetilde{P}\left(v_{1}\right), \ldots, \widetilde{P}\left(v_{N}\right)\right)
$$

is closed in $W^{k}(M, F)$. Take in particular $k=0$. Since $\mathscr{C}^{\infty}(M, F)$ is dense in $W^{d}(M, F)$, we see that in $W^{0}(M, F)$

$$
\left(\widetilde{P}\left(W^{d}(M, F)\right)\right)^{\perp}=\left(P\left(\mathscr{C}^{\infty}(M, F)\right)\right)^{\perp}=\operatorname{ker} \widetilde{P^{\star}}
$$

We have proved that

$$
\begin{equation*}
W^{0}(M, F)=\widetilde{P}\left(W^{d}(M, F)\right) \oplus \operatorname{ker} \widetilde{P^{\star}} \tag{2.6}
\end{equation*}
$$

Since $P^{\star}$ is also elliptic, it follows that ker $\widetilde{P^{\star}}$ is finite dimensional and that ker $\widetilde{P^{\star}}=$ ker $P^{\star}$ is contained in $\mathscr{C}^{\infty}(M, F)$. Thanks to Gårding's inequality, the decomposition formula (2.6) yields

$$
\begin{align*}
W^{k}(M, F) & =\widetilde{P}\left(W^{k+d}(M, F)\right) \oplus \operatorname{ker} P^{\star}  \tag{2.7}\\
\mathscr{C}^{\infty}(M, F) & =P\left(\mathscr{C}^{\infty}(M, F)\right) \oplus \operatorname{ker} P^{\star} \tag{2.8}
\end{align*}
$$

## § 3. Harmonic Forms and Hodge Theory on Riemannian Manifolds

## §3.1. Euclidean Structure of the Exterior Algebra

Let $(M, g)$ be an oriented Riemannian $\mathscr{C}^{\infty}$-manifold, $\operatorname{dim}_{\mathbb{R}} M=m$, and $E \rightarrow M$ a hermitian vector bundle of rank $r$ over $M$. We denote respectively by $\left(\xi_{1}, \ldots, \xi_{m}\right)$ and $\left(e_{1}, \ldots, e_{r}\right)$ orthonormal frames of $T_{M}$ and $E$ over an open subset $\Omega \subset M$, and by $\left(\xi_{1}^{\star}, \ldots, \xi_{m}^{\star}\right),\left(e_{1}^{\star}, \ldots, e_{r}^{\star}\right)$ the corresponding dual frames of $T_{M}^{\star}, E^{\star}$. Let $d V$ stand for the Riemannian volume form on $M$. The exterior algebra $\Lambda T_{M}^{\star}$ has a natural inner product $\langle\bullet, \bullet\rangle$ such that

$$
\begin{equation*}
\left\langle u_{1} \wedge \ldots \wedge u_{p}, v_{1} \wedge \ldots \wedge v_{p}\right\rangle=\operatorname{det}\left(\left\langle u_{j}, v_{k}\right\rangle\right)_{1 \leqslant j, k \leqslant p}, \quad u_{j}, v_{k} \in T_{M}^{\star} \tag{3.1}
\end{equation*}
$$

for all $p$, with $\Lambda T_{M}^{\star}=\bigoplus \Lambda^{p} T_{M}^{\star}$ as an orthogonal sum. Then the covectors $\xi_{I}^{\star}=\xi_{i_{1}}^{\star} \wedge$ $\cdots \wedge \xi_{i_{p}}^{\star}, i_{1}<i_{2}<\cdots<i_{p}$, provide an orthonormal basis of $\Lambda T_{M}^{\star}$. We also denote by $\langle\bullet, \bullet\rangle$ the corresponding inner product on $\Lambda T_{M}^{\star} \otimes E$.
(3.2) Hodge Star Operator. The Hodge-Poincaré-De Rham operator $\star$ is the collection of linear maps defined by

$$
\star: \Lambda^{p} T_{M}^{\star} \rightarrow \Lambda^{m-p} T_{M}^{\star}, \quad u \wedge \star v=\langle u, v\rangle d V, \quad \forall u, v \in \Lambda^{p} T_{M}^{\star}
$$

The existence and uniqueness of this operator is easily seen by using the duality pairing

$$
\begin{align*}
\Lambda^{p} T_{M}^{\star} \times \Lambda^{m-p} T_{M}^{\star} & \longrightarrow \mathbb{R} \\
(u, v) & \longmapsto u \wedge v / d V=\sum \varepsilon(I, \complement I) u_{I} v_{C I}, \tag{3.3}
\end{align*}
$$

where $u=\sum_{|I|=p} u_{I} \xi_{I}^{\star}, v=\sum_{|J|=m-p} v_{J} \xi_{J}^{\star}$, where $\lceil I$ stands for the (ordered) complementary multi-index of $I$ and $\varepsilon(I, \complement I)$ for the signature of the permutation $(1,2, \ldots, m) \mapsto$ ( $I, C I$ ). From this, we find

$$
\begin{equation*}
\star v=\sum_{|I|=p} \varepsilon(I, \complement I) v_{I} \xi_{\mathrm{C}_{I}}^{\star} \tag{3.4}
\end{equation*}
$$

More generally, the sesquilinear pairing $\{\bullet \bullet \bullet\}$ defined in (V-7.1) yields an operator $\star$ on vector valued forms, such that

$$
\left(3.3^{\prime}\right) \star: \Lambda^{p} T_{M}^{\star} \otimes E \rightarrow \Lambda^{m-p} T_{M}^{\star} \otimes E, \quad\{s, \star t\}=\langle s, t\rangle d V, \quad s, t \in \Lambda^{p} T_{M}^{\star} \otimes E,
$$

$$
\left(3.4^{\prime}\right) \star t=\sum_{|I|=p, \lambda} \varepsilon(I, \complement I) t_{I, \lambda} \xi_{\text {СI }}^{\star} \otimes e_{\lambda}
$$

for $t=\sum t_{I, \lambda} \xi_{I}^{\star} \otimes e_{\lambda}$. Since $\varepsilon(I, \complement I) \varepsilon(\complement I, I)=(-1)^{p(m-p)}=(-1)^{p(m-1)}$, we get immediately

$$
\begin{equation*}
\star \star t=(-1)^{p(m-1)} t \quad \text { on } \quad \Lambda^{p} T_{M}^{\star} \otimes E . \tag{3.5}
\end{equation*}
$$

It is clear that $\star$ is an isometry of $\Lambda^{\bullet} T_{M}^{\star} \otimes E$.
We shall also need a variant of the $\star$ operator, namely the conjugate-linear operator

$$
\#: \quad \Lambda^{p} T_{M}^{\star} \otimes E \longrightarrow \Lambda^{m-p} T_{M}^{\star} \otimes E^{\star}
$$

defined by $s \wedge \# t=\langle s, t\rangle d V$, where the wedge product $\wedge$ is combined with the canonical pairing $E \times E^{\star} \rightarrow \mathbb{C}$. We have

$$
\begin{equation*}
\# t=\sum_{|I|=p, \lambda} \varepsilon(I, C I) \bar{t}_{I, \lambda} \xi_{\mathrm{C}_{I}}^{\star} \otimes e_{\lambda}^{\star} . \tag{3.6}
\end{equation*}
$$

(3.7) Contraction by a Vector Field.. Given a tangent vector $\theta \in T_{M}$ and a form $u \in \Lambda^{p} T_{M}^{\star}$, the contraction $\left.\theta\right\lrcorner u \in \Lambda^{p-1} T_{M}^{\star}$ is defined by

$$
\theta\lrcorner u\left(\eta_{1}, \ldots, \eta_{p-1}\right)=u\left(\theta, \eta_{1}, \ldots, \eta_{p-1}\right), \quad \eta_{j} \in T_{M} .
$$

In terms of the basis $\left(\xi_{j}\right), \bullet \bullet$ is the bilinear operation characterized by

$$
\left.\xi_{l}\right\lrcorner\left(\xi_{i_{1}}^{\star} \wedge \ldots \wedge \xi_{i_{p}}^{\star}\right)= \begin{cases}0 & \text { if } l \notin\left\{i_{1}, \ldots, i_{p}\right\} \\ (-1)^{k-1} \xi_{i_{1}}^{\star} \wedge \ldots \widehat{\xi_{i_{k}}^{\star}} \ldots \wedge \xi_{i_{p}}^{\star} & \text { if } l=i_{k} .\end{cases}
$$

This formula is in fact valid even when $\left(\xi_{j}\right)$ is non orthonormal. A rather easy computation shows that $\theta\lrcorner \bullet$ is a derivation of the exterior algebra, i.e. that

$$
\left.\theta\lrcorner(u \wedge v)=(\theta\lrcorner u) \wedge v+(-1)^{\operatorname{deg} u} u \wedge(\theta\lrcorner v\right) .
$$

Moreover, if $\widetilde{\theta}=\langle\bullet, \theta\rangle \in T_{M}^{\star}$, the operator $\left.\theta\right\lrcorner \bullet$ is the adjoint map of $\widetilde{\theta} \wedge \bullet$, that is,

$$
\begin{equation*}
\langle\theta\lrcorner u, v\rangle=\langle u, \widetilde{\theta} \wedge v\rangle, \quad u, v \in \Lambda T_{M}^{\star} \tag{3.8}
\end{equation*}
$$

Indeed, this property is immediately checked when $\theta=\xi_{l}, u=\xi_{I}^{\star}, v=\xi_{J}^{\star}$.

## § 3.2. Laplace-Beltrami Operators

Let us consider the Hilbert space $L^{2}\left(M, \Lambda^{p} T_{M}^{\star}\right)$ of $p$-forms $u$ on $M$ with measurable coefficients such that

$$
\|u\|^{2}=\int_{M}|u|^{2} d V<+\infty .
$$

We denote by $\langle\langle\rangle$,$\rangle the global inner product on L^{2}$-forms. One can define similarly the Hilbert space $L^{2}\left(M, \Lambda^{p} T_{M}^{\star} \otimes E\right)$.
(3.9) Theorem. The operator $d^{\star}=(-1)^{m p+1} \star d \star$ is the formal adjoint of the exterior derivative $d$ acting on $\mathscr{C}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star} \otimes E\right)$.

Proof. If $u \in \mathscr{C}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star}\right), v \in \mathscr{C}^{\infty}\left(M, \Lambda^{p+1} T_{M}^{\star} \otimes\right)$ are compactly supported we get

$$
\begin{aligned}
\langle\langle d u, v\rangle\rangle & =\int_{M}\langle d u, v\rangle d V=\int_{M} d u \wedge \star v \\
& =\int_{M} d(u \wedge \star v)-(-1)^{p} u \wedge d \star v=-(-1)^{p} \int_{M} u \wedge d \star v
\end{aligned}
$$

by Stokes' formula. Therefore (3.4) implies

$$
\langle\langle d u, v\rangle\rangle=-(-1)^{p}(-1)^{p(m-1)} \int_{M} u \wedge \star \star d \star v=(-1)^{m p+1}\langle\langle u, \star d \star v\rangle\rangle
$$

(3.10) Remark. If $m$ is even, the formula reduces to $d^{\star}=-\star d \star$.
(3.11) Definition. The operator $\Delta=d d^{\star}+d^{\star} d$ is called the Laplace-Beltrami operator of $M$.

Since $d^{\star}$ is the adjoint of $d$, the Laplace operator $\Delta$ is formally self-adjoint, i.e. $\langle\langle\Delta u, v\rangle\rangle=\langle\langle u, \Delta v\rangle\rangle$ when the forms $u, v$ are of class $\mathscr{C}^{\infty}$ and compactly supported.
(3.12) Example. Let $M$ be an open subset of $\mathbb{R}^{m}$ and $g=\sum_{i=1}^{m} d x_{i}^{2}$. In that case we get

$$
\begin{aligned}
u & =\sum_{|I|=p} u_{I} d x_{I}, \quad d u=\sum_{|I|=p, j} \frac{\partial u_{I}}{\partial x_{j}} d x_{j} \wedge d x_{I}, \\
\langle\langle u, v\rangle\rangle & =\int_{M}\langle u, v\rangle d V=\int_{M} \sum_{I} u_{I} v_{I} d V
\end{aligned}
$$

One can write $d v=\sum d x_{j} \wedge\left(\partial v / \partial x_{j}\right)$ where $\partial v / \partial x_{j}$ denotes the form $v$ in which all coefficients $v_{I}$ are differentiated as $\partial v_{I} / \partial x_{j}$. An integration by parts combined with contraction gives

$$
\begin{aligned}
\left\langle\left\langle d^{\star} u, v\right\rangle\right\rangle & =\langle\langle u, d v\rangle\rangle=\int_{M}\left\langle u, \sum_{j} d x_{j} \wedge \frac{\partial v}{\partial x_{j}}\right\rangle d V \\
& \left.\left.=\int_{M} \sum_{j}\left\langle\frac{\partial}{\partial x_{j}}\right\lrcorner u, \frac{\partial v}{\partial x_{j}}\right\rangle d V=-\int_{M}\left\langle\sum_{j} \frac{\partial}{\partial x_{j}}\right\lrcorner \frac{\partial u}{\partial x_{j}}, v\right\rangle d V \\
d^{\star} u & \left.\left.=-\sum_{j} \frac{\partial}{\partial x_{j}}\right\lrcorner \frac{\partial u}{\partial x_{j}}=-\sum_{I, j} \frac{\partial u_{I}}{\partial x_{j}} \frac{\partial}{\partial x_{j}}\right\lrcorner d x_{I} .
\end{aligned}
$$

We get therefore

$$
\begin{aligned}
d d^{\star} u & \left.=-\sum_{I, j, k} \frac{\partial^{2} u_{I}}{\partial x_{j} \partial x_{k}} d x_{k} \wedge\left(\frac{\partial}{\partial x_{j}}\right\lrcorner d x_{I}\right), \\
d^{\star} d u & \left.=-\sum_{I, j, k} \frac{\partial^{2} u_{I}}{\partial x_{j} \partial x_{k}} \frac{\partial}{\partial x_{j}}\right\lrcorner\left(d x_{k} \wedge d x_{I}\right) .
\end{aligned}
$$

Since

$$
\left.\left.\left.\frac{\partial}{\partial x_{j}}\right\lrcorner\left(d x_{k} \wedge d x_{I}\right)=\left(\frac{\partial}{\partial x_{j}}\right\lrcorner d x_{k}\right) d x_{I}-d x_{k} \wedge\left(\frac{\partial}{\partial x_{j}}\right\lrcorner d x_{I}\right)
$$

we obtain

$$
\Delta u=-\sum_{I}\left(\sum_{j} \frac{\partial^{2} u_{I}}{\partial x_{j}^{2}}\right) d x_{I}
$$

In the case of an arbitrary riemannian manifold $(M, g)$ we have

$$
\begin{aligned}
u & =\sum u_{I} \xi_{I}^{\star}, \\
d u & =\sum_{I, j}\left(\xi_{j} \cdot u_{I}\right) \xi_{j}^{\star} \wedge \xi_{I}^{\star}+\sum_{I} u_{I} d \xi_{I}^{\star}, \\
d^{\star} u & \left.=-\sum_{I, j}\left(\xi_{j} \cdot u_{I}\right) \xi_{j}\right\lrcorner \xi_{I}^{\star}+\sum_{I, K} \alpha_{I, K} u_{I} \xi_{K}^{\star},
\end{aligned}
$$

for some $\mathscr{C}^{\infty}$ coefficients $\alpha_{I, K},|I|=p,|K|=p-1$. It follows that the principal part of $\Delta$ is the same as that of the second order operator

$$
u \longmapsto-\sum_{I}\left(\sum_{j} \xi_{j}^{2} \cdot u_{I}\right) \xi_{I}^{\star} .
$$

As a consequence, $\Delta$ is elliptic.

Assume now that $D_{E}$ is a hermitian connection on $E$. The formal adjoint operator of $D_{E}$ acting on $\mathscr{C}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star} \otimes E\right)$ is

$$
\begin{equation*}
D_{E}^{\star}=(-1)^{m p+1} \star D_{E} \star . \tag{3.13}
\end{equation*}
$$

Indeed, if $s \in \mathscr{C}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star} \otimes E\right), t \in \mathscr{C}^{\infty}\left(M, \Lambda^{p+1} T_{M}^{\star} \otimes E\right)$ have compact support, we get

$$
\begin{aligned}
\left\langle\left\langle D_{E} s, t\right\rangle\right\rangle & =\int_{M}\left\langle D_{E} s, t\right\rangle d V=\int_{M}\left\{D_{E} s, \star t\right\} \\
& =\int_{M} d\{s, \star t\}-(-1)^{p}\left\{s, D_{E} \star t\right\}=(-1)^{m p+1}\left\langle\left\langle s, \star D_{E} \star t\right\rangle\right\rangle .
\end{aligned}
$$

(3.14) Definition. The Laplace-Beltrami operator associated to $D_{E}$ is the second order operator $\Delta_{E}=D_{E} D_{E}^{\star}+D_{E}^{\star} D_{E}$.
$\Delta_{E}$ is a self-adjoint elliptic operator with principal part

$$
s \longmapsto-\sum_{I, \lambda}\left(\sum_{j} \xi_{j}^{2} \cdot s_{I, \lambda}\right) \xi_{I}^{\star} \otimes e_{\lambda} .
$$

## $\S$ 3.3. Harmonic Forms and Hodge Isomorphism

Let $E$ be a hermitian vector bundle over a compact Riemannian manifold $(M, g)$. We assume that $E$ possesses a flat hermitian connection $D_{E}$ (this means that $\Theta\left(D_{E}\right)=$ $D_{E}^{2}=0$, or equivalently, that $E$ is given by a representation $\pi_{1}(M) \rightarrow U(r)$, cf. § V-6). A fundamental example is of course the trivial bundle $E=M \times \mathbb{C}$ with the connection $D_{E}=d$. Thanks to our flatness assumption, $D_{E}$ defines a generalized De Rham complex

$$
D_{E}: \mathscr{C}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star} \otimes E\right) \longrightarrow \mathscr{C}^{\infty}\left(M, \Lambda^{p+1} T_{M}^{\star} \otimes E\right) .
$$

The cohomology groups of this complex will be denoted by $H_{D R}^{p}(M, E)$.

The space of harmonic forms of degree $p$ with respect to the Laplace-Beltrami operator $\Delta_{E}=D_{E} D_{E}^{\star}+D_{E}^{\star} D_{E}$ is defined by

$$
\begin{equation*}
\mathscr{H}^{p}(M, E)=\left\{s \in \mathscr{C}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star} \otimes E\right) ; \Delta_{E} s=0\right\} . \tag{3.15}
\end{equation*}
$$

Since $\left\langle\left\langle\Delta_{E} s, s\right\rangle\right\rangle=\left\|D_{E} s\right\|^{2}+\left\|D_{E}^{\star} s\right\|^{2}$, we see that $s \in \mathscr{H}^{p}(M, E)$ if and only if $D_{E} s=$ $D_{E}^{\star} s=0$.
(3.16) Theorem. For any $p$, there exists an orthogonal decomposition

$$
\begin{aligned}
& \mathscr{C}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star} \otimes E\right)=\mathscr{H}^{p}(M, E) \oplus \operatorname{Im} D_{E} \oplus \operatorname{Im} D_{E}^{\star} \\
& \operatorname{Im} D_{E}=D_{E}\left(\mathscr{C}^{\infty}\left(M, \Lambda^{p-1} T_{M}^{\star} \otimes E\right)\right), \\
& \operatorname{Im} D_{E}^{\star}=D_{E}^{\star}\left(\mathscr{C}^{\infty}\left(M, \Lambda^{p+1} T_{M}^{\star} \otimes E\right)\right) .
\end{aligned}
$$

Proof. It is immediate that $\mathscr{H}^{p}(M, E)$ is orthogonal to both subspaces $\operatorname{Im} D_{E}$ and $\operatorname{Im} D_{E}^{\star}$. The orthogonality of these two subspaces is also clear, thanks to the assumption $D_{E}^{2}=0$ :

$$
\left\langle\left\langle D_{E} s, D_{E}^{\star} t\right\rangle\right\rangle=\left\langle\left\langle D_{E}^{2} s, t\right\rangle\right\rangle=0 .
$$

We apply now Cor. 2.4 to the elliptic operator $\Delta_{E}=\Delta_{E}^{\star}$ acting on $p$-forms, i.e. on the bundle $F=\Lambda^{p} T_{M}^{\star} \otimes E$. We get

$$
\begin{aligned}
& \mathscr{C}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star} \otimes E\right)=\mathscr{H}^{p}(M, E) \oplus \Delta_{E}\left(\mathscr{C}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star} \otimes E\right)\right) \\
& \operatorname{Im} \Delta_{E}=\operatorname{Im}\left(D_{E} D_{E}^{\star}+D_{E}^{\star} D_{E}\right) \subset \operatorname{Im} D_{E}+\operatorname{Im} D_{E}^{\star}
\end{aligned}
$$

(3.17) Hodge isomorphism theorem. The De Rham cohomology group $H_{D R}^{p}(M, E)$ is finite dimensional and $H_{D R}^{p}(M, E) \simeq \mathscr{H}^{p}(M, E)$.

Proof. According to decomposition 3.16, we get

$$
\begin{aligned}
& B_{D R}^{p}(M, E)=D_{E}\left(\mathscr{C}^{\infty}\left(M, \Lambda^{p-1} T_{M}^{\star} \otimes E\right)\right) \\
& Z_{D R}^{p}(M, E)=\operatorname{ker} D_{E}=\left(\operatorname{Im} D_{E}^{\star}\right)^{\perp}=\mathscr{H}^{p}(M, E) \oplus \operatorname{Im} D_{E}
\end{aligned}
$$

This shows that every De Rham cohomology class contains a unique harmonic representative.
(3.18) Poincaré duality. The bilinear pairing

$$
H_{D R}^{p}(M, E) \times H_{D R}^{m-p}\left(M, E^{\star}\right) \longrightarrow \mathbb{C}, \quad(s, t) \longmapsto \int_{M} s \wedge t
$$

is a non degenerate duality.
Proof. First note that there exists a naturally defined flat connection $D_{E^{\star}}$ such that for any $s_{1} \in \mathscr{C}_{\bullet}^{\infty}(M, E), s_{2} \in \mathscr{C}_{\bullet}^{\infty}\left(M, E^{\star}\right)$ we have

$$
\begin{equation*}
d\left(s_{1} \wedge s_{2}\right)=\left(D_{E} s_{1}\right) \wedge s_{2}+(-1)^{\operatorname{deg} s_{1}} s_{1} \wedge D_{E^{\star}} s_{2} \tag{3.19}
\end{equation*}
$$

It is then a consequence of Stokes' formula that the map $(s, t) \mapsto \int_{M} s \wedge t$ can be factorized through cohomology groups. Let $s \in \mathscr{C}^{\infty}\left(M, \Lambda^{p} T_{M}^{\star} \otimes E\right)$. We leave to the reader the proof of the following formulas (use (3.19) in analogy with the proof of Th. 3.9):

$$
\begin{align*}
D_{E^{\star}}(\# s) & =(-1)^{p} \# D_{E}^{\star} s, \\
\delta_{E^{\star}}(\# s) & =(-1)^{p+1} \# D_{E} s,  \tag{3.20}\\
\Delta_{E^{\star}}(\# s) & =\# \Delta_{E} s,
\end{align*}
$$

Consequently $\# s \in \mathscr{H}^{m-p}\left(M, E^{\star}\right)$ if and only if $s \in \mathscr{H}^{p}(M, E)$. Since

$$
\int_{M} s \wedge \# s=\int_{M}|s|^{2} d V=\|s\|^{2}
$$

we see that the Poincaré pairing has zero kernel in the left hand factor $\mathscr{H}^{p}(M, E) \simeq$ $H_{D R}^{p}(M, E)$. By symmetry, it has also zero kernel on the right. The proof is achieved.

## § 4. Hermitian and Kähler Manifolds

Let $X$ be a complex $n$-dimensional manifold. A hermitian metric on $X$ is a positive definite hermitian form of class $\mathscr{C}^{\infty}$ on $T_{X}$; in a coordinate system $\left(z_{1}, \ldots, z_{n}\right)$, such a form can be written $h(z)=\sum_{1 \leqslant j, k \leqslant n} h_{j k}(z) d z_{j} \otimes d \bar{z}_{k}$, where $\left(h_{j k}\right)$ is a positive hermitian matrix with $\mathscr{C}^{\infty}$ coefficients. According to (III-1.8), the fundamental ( 1,1 )-form associated to $h$ is the positive form of type $(1,1)$

$$
\omega=-\operatorname{Im} h=\frac{\mathrm{i}}{2} \sum h_{j k} d z_{j} \wedge d \bar{z}_{k}, \quad 1 \leqslant j, k \leqslant n
$$

## (4.1) Definition.

a) A hermitian manifold is a pair $(X, \omega)$ where $\omega$ is a $\mathscr{C}^{\infty}$ positive definite $(1,1)$-form on $X$.
b) The metric $\omega$ is said to be kähler if $d \omega=0$.
c) $X$ is said to be a Kähler manifold if $X$ carries at least one Kähler metric.

Since $\omega$ is real, the conditions $d \omega=0, d^{\prime} \omega=0, d^{\prime \prime} \omega=0$ are all equivalent. In local coordinates we see that $d^{\prime} \omega=0$ if and only if

$$
\frac{\partial h_{j k}}{\partial z_{l}}=\frac{\partial h_{l k}}{\partial z_{j}} \quad, \quad 1 \leqslant j, k, l \leqslant n
$$

A simple computation gives

$$
\frac{\omega^{n}}{n!}=\operatorname{det}\left(h_{j k}\right) \bigwedge_{1 \leqslant j \leqslant n}\left(\frac{\mathrm{i}}{2} d z_{j} \wedge d \bar{z}_{j}\right)=\operatorname{det}\left(h_{j k}\right) d x_{1} \wedge d y_{1} \wedge \cdots \wedge d x_{n} \wedge d y_{n}
$$

where $z_{n}=x_{n}+\mathrm{i} y_{n}$. Therefore the $(n, n)$-form

$$
\begin{equation*}
d V=\frac{1}{n!} \omega^{n} \tag{4.2}
\end{equation*}
$$

is positive and coincides with the hermitian volume element of $X$. If $X$ is compact, then $\int_{X} \omega^{n}=n!\operatorname{Vol}_{\omega}(X)>0$. This simple remark already implies that compact Kähler manifolds must satisfy some restrictive topological conditions:

## (4.3) Consequence.

a) If $(X, \omega)$ is compact Kähler and if $\{\omega\}$ denotes the cohomology class of $\omega$ in $H^{2}(X, \mathbb{R})$, then $\{\omega\}^{n} \neq 0$.
b) If $X$ is compact Kähler, then $H^{2 k}(X, \mathbb{R}) \neq 0$ for $0 \leqslant k \leqslant n$. In fact, $\{\omega\}^{k}$ is a non zero class in $H^{2 k}(X, \mathbb{R})$.
(4.4) Example. The complex projective space $\mathbb{P}^{n}$ is Kähler. A natural Kähler metric $\omega$ on $\mathbb{P}^{n}$, called the Fubini-Study metric, is defined by

$$
p^{\star} \omega=\frac{\mathrm{i}}{2 \pi} d^{\prime} d^{\prime \prime} \log \left(\left|\zeta_{0}\right|^{2}+\left|\zeta_{1}\right|^{2}+\cdots+\left|\zeta_{n}\right|^{2}\right)
$$

where $\zeta_{0}, \zeta_{1}, \ldots, \zeta_{n}$ are coordinates of $\mathbb{C}^{n+1}$ and where $p: \mathbb{C}^{n+1} \backslash\{0\} \rightarrow \mathbb{P}^{n}$ is the projection. Let $z=\left(\zeta_{1} / \zeta_{0}, \ldots, \zeta_{n} / \zeta_{0}\right)$ be non homogeneous coordinates on $\mathbb{C}^{n} \subset \mathbb{P}^{n}$. Then (V-15.8) and (V-15.12) show that

$$
\omega=\frac{\mathrm{i}}{2 \pi} d^{\prime} d^{\prime \prime} \log \left(1+|z|^{2}\right)=\frac{\mathrm{i}}{2 \pi} c(\mathbb{G}(1)), \quad \int_{\mathbb{P}^{n}} \omega^{n}=1 .
$$

Furthermore $\{\omega\} \in H^{2}\left(\mathbb{P}^{n}, \mathbb{Z}\right)$ is a generator of the cohomology algebra $H^{\bullet}\left(\mathbb{P}^{n}, \mathbb{Z}\right)$ in virtue of Th. V-15.10.
(4.5) Example. A complex torus is a quotient $X=\mathbb{C}^{n} / \Gamma$ by a lattice $\Gamma$ of rank $2 n$. Then $X$ is a compact complex manifold. Any positive definite hermitian form $\omega=$ i $\sum h_{j k} d z_{j} \wedge d \bar{z}_{k}$ with constant coefficients defines a Kähler metric on $X$.
(4.6) Example. Every (complex) submanifold $Y$ of a Kähler manifold $(X, \omega)$ is Kähler with metric $\omega_{\mid Y}$. Especially, all submanifolds of $\mathbb{P}^{n}$ are Kähler.
(4.7) Example. Consider the complex surface

$$
X=\left(\mathbb{C}^{2} \backslash\{0\}\right) / \Gamma
$$

where $\Gamma=\left\{\lambda^{n} ; n \in \mathbb{Z}\right\}, \lambda<1$, acts as a group of homotheties. Since $\mathbb{C}^{2} \backslash\{0\}$ is diffeomorphic to $\mathbb{R}_{+}^{\star} \times S^{3}$, we have $X \simeq S^{1} \times S^{3}$. Therefore $H^{2}(X, \mathbb{R})=0$ by Künneth's formula IV-15.10, and property 4.3 b ) shows that $X$ is not Kähler. More generally, one can obtaintake $\Gamma$ to be an infinite cyclic group generated by a holomorphic contraction of $\mathbb{C}^{2}$, of the form

$$
\binom{z_{1}}{z_{2}} \longmapsto\binom{\lambda_{1} z_{1}}{\lambda_{2} z_{2}}, \quad \text { resp. } \quad\binom{z_{1}}{z_{2}} \longmapsto\binom{\lambda z_{1}}{\lambda z_{2}+z_{1}^{p}},
$$

where $\lambda, \lambda_{1}, \lambda_{2}$ are complex numbers such that $0<\left|\lambda_{1}\right| \leqslant\left|\lambda_{2}\right|<1,0<|\lambda|<1$, and $p$ a positive integer. These non Kähler surfaces are called Hopf surfaces.

The following Theorem shows that a hermitian metric $\omega$ on $X$ is Kähler if and only if the metric $\omega$ is tangent at order 2 to a hermitian metric with constant coefficients at every point of $X$.
(4.8) Theorem. Let $\omega$ be a $\mathscr{C}^{\infty}$ positive definite $(1,1)$-form on $X$. In order that $\omega$ be Kähler, it is necessary and sufficient that to every point $x_{0} \in X$ corresponds a holomorphic coordinate system $\left(z_{1}, \ldots, z_{n}\right)$ centered at $x_{0}$ such that

$$
\begin{equation*}
\omega=\mathrm{i} \sum_{1 \leqslant l, m \leqslant n} \omega_{l m} d z_{l} \wedge d \bar{z}_{m}, \quad \omega_{l m}=\delta_{l m}+O\left(|z|^{2}\right) . \tag{4.9}
\end{equation*}
$$

If $\omega$ is Kähler, the coordinates $\left(z_{j}\right)_{1 \leqslant j \leqslant n}$ can be chosen such that

$$
\begin{equation*}
\omega_{l m}=\left\langle\frac{\partial}{\partial z_{l}}, \frac{\partial}{\partial z_{m}}\right\rangle=\delta_{l m}-\sum_{1 \leqslant j, k \leqslant n} c_{j k l m} z_{j} \bar{z}_{k}+O\left(|z|^{3}\right) \tag{4.10}
\end{equation*}
$$

where ( $c_{j k l m}$ ) are the coefficients of the Chern curvature tensor

$$
\begin{equation*}
\Theta\left(T_{X}\right)_{x_{0}}=\sum_{j, k, l, m} c_{j k l m} d z_{j} \wedge d \bar{z}_{k} \otimes\left(\frac{\partial}{\partial z_{l}}\right)^{\star} \otimes \frac{\partial}{\partial z_{m}} \tag{4.11}
\end{equation*}
$$

associated to $\left(T_{X}, \omega\right)$ at $x_{0}$. Such a system $\left(z_{j}\right)$ will be called a geodesic coordinate system at $x_{0}$.

Proof. It is clear that (4.9) implies $d_{x_{0}} \omega=0$, so the condition is sufficient. Assume now that $\omega$ is Kähler. Then one can choose local coordinates $\left(x_{1}, \ldots, x_{n}\right)$ such that $\left(d x_{1}, \ldots, d x_{n}\right)$ is an $\omega$-orthonormal basis of $T_{x_{0}}^{\star} X$. Therefore

$$
\begin{align*}
\omega & =\mathrm{i} \sum_{1 \leqslant l, m \leqslant n} \widetilde{\omega}_{l m} d x_{l} \wedge d \bar{x}_{m}, \quad \text { where } \\
\widetilde{\omega}_{l m} & =\delta_{l m}+O(|x|)=\delta_{l m}+\sum_{1 \leqslant j \leqslant n}\left(a_{j l m} x_{j}+a_{j l m}^{\prime} \bar{x}_{j}\right)+O\left(|x|^{2}\right) . \tag{4.12}
\end{align*}
$$

Since $\omega$ is real, we have $a_{j l m}^{\prime}=\bar{a}_{j m l}$; on the other hand the Kähler condition $\partial \omega_{l m} / \partial x_{j}=$ $\partial \omega_{j m} / \partial x_{l}$ at $x_{0}$ implies $a_{j l m}=a_{l j m}$. Set now

$$
z_{m}=x_{m}+\frac{1}{2} \sum_{j, l} a_{j l m} x_{j} x_{l}, \quad 1 \leqslant m \leqslant n .
$$

Then $\left(z_{m}\right)$ is a coordinate system at $x_{0}$, and

$$
\begin{aligned}
d z_{m}= & d x_{m}+\sum_{j, l} a_{j l m} x_{j} d x_{l}, \\
\mathrm{i} \sum_{m} d z_{m} \wedge d \bar{z}_{m}= & \mathrm{i} \sum_{m} d x_{m} \wedge d \bar{x}_{m}+\mathrm{i} \sum_{j, l, m} a_{j l m} x_{j} d x_{l} \wedge d \bar{x}_{m} \\
& +\mathrm{i} \sum_{j, l, m} \bar{a}_{j l m} \bar{x}_{j} d x_{m} \wedge d \bar{x}_{l}+O\left(|x|^{2}\right) \\
= & \mathrm{i} \sum_{l, m} \widetilde{\omega}_{l m} d x_{l} \wedge \bar{d} x_{m}+O\left(|x|^{2}\right)=\omega+O\left(|z|^{2}\right)
\end{aligned}
$$

Condition (4.9) is proved. Suppose the coordinates $\left(x_{m}\right)$ chosen from the beginning so that (4.9) holds with respect to $\left(x_{m}\right)$. Then the Taylor expansion (4.12) can be refined into

$$
\begin{align*}
\widetilde{\omega}_{l m} & =\delta_{l m}+O\left(|x|^{2}\right)  \tag{4.13}\\
& =\delta_{l m}+\sum_{j, k}\left(a_{j k l m} x_{j} \bar{x}_{k}+a_{j k l m}^{\prime} x_{j} x_{k}+a_{j k l m}^{\prime \prime} \bar{x}_{j} \bar{x}_{k}\right)+O\left(|x|^{3}\right) .
\end{align*}
$$

These new coefficients satisfy the relations

$$
a_{j k l m}^{\prime}=a_{k j l m}^{\prime}, \quad a_{j k l m}^{\prime \prime}=\bar{a}_{j k m l}^{\prime}, \quad \bar{a}_{j k l m}=a_{k j m l} .
$$

The Kähler condition $\partial \omega_{l m} / \partial x_{j}=\partial \omega_{j m} / \partial x_{l}$ at $x=0$ gives the equality $a_{j k l m}^{\prime}=a_{l k j m}^{\prime}$; in particular $a_{j k l m}^{\prime}$ is invariant under all permutations of $j, k, l$. If we set

$$
z_{m}=x_{m}+\frac{1}{3} \sum_{j, k, l} a_{j k l m}^{\prime} x_{j} x_{k} x_{l}, \quad 1 \leqslant m \leqslant n
$$

then by (4.13) we find

$$
\begin{align*}
d z_{m} & =d x_{m}+\sum_{j, k, l} a_{j k l m}^{\prime} x_{j} x_{k} d x_{l}, \quad 1 \leqslant m \leqslant n, \\
\omega & =\mathrm{i} \sum_{1 \leqslant m \leqslant n} d z_{m} \wedge d \bar{z}_{m}+\mathrm{i} \sum_{j, k, l, m} a_{j k l m} x_{j} \bar{x}_{k} d x_{l} \wedge d \bar{x}_{m}+O\left(|x|^{3}\right), \\
\omega & =\mathrm{i} \sum_{1 \leqslant m \leqslant n} d z_{m} \wedge d \bar{z}_{m}+\mathrm{i} \sum_{j, k, l, m} a_{j k l m} z_{j} \bar{z}_{k} d z_{l} \wedge d \bar{z}_{m}+O\left(|z|^{3}\right) . \tag{4.14}
\end{align*}
$$

It is now easy to compute the Chern curvature tensor $\Theta\left(T_{X}\right)_{x_{0}}$ in terms of the coefficients $a_{j k l m}$. Indeed

$$
\begin{aligned}
\left\langle\frac{\partial}{\partial z_{l}}, \frac{\partial}{\partial z_{m}}\right\rangle & =\delta_{l m}+\sum_{j, k} a_{j k l m} z_{j} \bar{z}_{k}+O\left(|z|^{3}\right), \\
d^{\prime}\left\langle\frac{\partial}{\partial z_{l}}, \frac{\partial}{\partial z_{m}}\right\rangle & =\left\{D^{\prime} \frac{\partial}{\partial z_{l}}, \frac{\partial}{\partial z_{m}}\right\}=\sum_{j, k} a_{j k l m} \bar{z}_{k} d z_{j}+O\left(|z|^{2}\right), \\
\Theta\left(T_{X}\right) \cdot \frac{\partial}{\partial z_{l}} & =D^{\prime \prime} D^{\prime}\left(\frac{\partial}{\partial z_{l}}\right)=-\sum_{j, k, m} a_{j k l m} d z_{j} \wedge d \bar{z}_{k} \otimes \frac{\partial}{\partial z_{m}}+O(|z|),
\end{aligned}
$$

therefore $c_{j k l m}=-a_{j k l m}$ and the expansion (4.10) follows from (4.14).
(4.15) Remark. As a by-product of our computations, we find that on a Kähler manifold the coefficients of $\Theta\left(T_{X}\right)$ satisfy the symmetry relations

$$
\bar{c}_{j k l m}=c_{k j m l}, \quad c_{j k l m}=c_{l k j m}=c_{j m l k}=c_{l m j k}
$$

## § 5. Basic Results of Kähler Geometry

## § 5.1. Operators of Hermitian Geometry

Let $(X, \omega)$ be a hermitian manifold and let $z_{j}=x_{j}+\mathrm{i} y_{j}, 1 \leqslant j \leqslant n$, be analytic coordinates at a point $x \in X$ such that $\omega(x)=\mathrm{i} \sum d z_{j} \wedge d \bar{z}_{j}$ is diagonalized at this point. The associated hermitian form is the $h(x)=2 \sum d z_{j} \otimes d \bar{z}_{j}$ and its real part is the euclidean metric $2 \sum\left(d x_{j}\right)^{2}+\left(d y_{j}\right)^{2}$. It follows from this that $\left|d x_{j}\right|=\left|d y_{j}\right|=1 / \sqrt{2}$, $\left|d z_{j}\right|=\left|d \bar{z}_{j}\right|=1$, and that $\left(\partial / \partial z_{1}, \ldots, \partial / \partial z_{n}\right)$ is an orthonormal basis of $\left(T_{x}^{\star} X, \omega\right)$. Formula (3.1) with $u_{j}, v_{k}$ in the orthogonal sum $\left(\mathbb{C} \otimes T_{X}\right)^{\star}=T_{X}^{\star} \oplus \overline{T_{X}^{\star}}$ defines a natural inner product on the exterior algebra $\Lambda^{\bullet}\left(\mathbb{C} \otimes T_{X}\right)^{\star}$. The norm of a form

$$
u=\sum_{I, J} u_{I, J} d z_{I} \wedge d \bar{z}_{J} \in \Lambda\left(\mathbb{C} \otimes T_{X}\right)^{\star}
$$

at the given point $x$ is then equal to

$$
\begin{equation*}
|u(x)|^{2}=\sum_{I, J}\left|u_{I, J}(x)\right|^{2} \tag{5.1}
\end{equation*}
$$

The Hodge $\star$ operator (3.2) can be extended to $\mathbb{C}$-valued forms by the formula

$$
\begin{equation*}
u \wedge \star \bar{v}=\langle u, v\rangle d V . \tag{5.2}
\end{equation*}
$$

It follows that $\star$ is a $\mathbb{C}$-linear isometry

$$
\star: \quad \Lambda^{p, q} T_{X}^{\star} \longrightarrow \Lambda^{n-q, n-p} T_{X}^{\star} .
$$

The usual operators of hermitian geometry are the operators $d, \delta=-\star d \star, \Delta=d \delta+\delta d$ already defined, and their complex counterparts

$$
\left\{\begin{align*}
d & =d^{\prime}+d^{\prime \prime},  \tag{5.3}\\
\delta & =d^{\prime \star}+d^{\prime \prime \star}, \quad d^{\prime \star}=\left(d^{\prime}\right)^{\star}=-\star d^{\prime \prime} \star, \quad d^{\prime \prime \star}=\left(d^{\prime \prime}\right)^{\star}=-\star d^{\prime} \star, \\
\Delta^{\prime} & =d^{\prime} d^{\prime \star}+d^{\prime \star} d^{\prime}, \quad \Delta^{\prime \prime}=d^{\prime \prime} d^{\prime \prime \star}+d^{\prime \prime \star} d^{\prime \prime} .
\end{align*}\right.
$$

Another important operator is the operator $L$ of type $(1,1)$ defined by

$$
\begin{equation*}
L u=\omega \wedge u \tag{5.4}
\end{equation*}
$$

and its adjoint $\Lambda=\star^{-1} L \star$ :

$$
\begin{equation*}
\langle u, \Lambda v\rangle=\langle L u, v\rangle . \tag{5.5}
\end{equation*}
$$

## §5.2. Commutation Identities

If $A, B$ are endomorphisms of the algebra $\mathscr{C}_{\bullet, \bullet}^{\infty}(X, \mathbb{C})$, their graded commutator (or graded Lie bracket) is defined by

$$
\begin{equation*}
[A, B]=A B-(-1)^{a b} B A \tag{5.6}
\end{equation*}
$$

where $a, b$ are the degrees of $A$ and $B$ respectively. If $C$ is another endomorphism of degree $c$, the following Jacobi identity is easy to check:

$$
\begin{equation*}
(-1)^{c a}[A,[B, C]]+(-1)^{a b}[B,[C, A]]+(-1)^{b c}[C,[A, B]]=0 \tag{5.7}
\end{equation*}
$$

For any $\alpha \in \Lambda^{p, q} T_{X}^{\star}$, we still denote by $\alpha$ the endomorphism of type $(p, q)$ on $\Lambda^{\bullet \bullet \bullet} T_{X}^{\star}$ defined by $u \mapsto \alpha \wedge u$. Consider further a real (1,1)-form

$$
\gamma \in \Lambda^{1,1} T_{X}^{\star} .
$$

There exists an $\omega$-orthogonal basis $\left(\zeta_{1}, \zeta_{2}, \ldots, \zeta_{n}\right)$ in $T_{X}$ which diagonalizes both forms $\omega$ and $\gamma$ :

$$
\omega=\mathrm{i} \sum_{1 \leqslant j \leqslant n} \zeta_{j}^{\star} \wedge \bar{\zeta}_{j}^{\star}, \quad \gamma=\mathrm{i} \sum_{1 \leqslant j \leqslant n} \gamma_{j} \zeta_{j}^{\star} \wedge \bar{\zeta}_{j}^{\star}, \quad \gamma_{j} \in \mathbb{R} .
$$

(5.8) Proposition. For every form $u=\sum u_{J, K} \zeta_{J}^{\star} \wedge \bar{\zeta}_{K}^{\star}$, one has

$$
[\gamma, \Lambda] u=\sum_{J, K}\left(\sum_{j \in J} \gamma_{j}+\sum_{j \in K} \gamma_{j}-\sum_{1 \leqslant j \leqslant n} \gamma_{j}\right) u_{J, K} \zeta_{J}^{\star} \wedge \bar{\zeta}_{K}^{\star} .
$$

Proof. If $u$ is of type $(p, q)$, a brute-force computation yields

$$
\begin{aligned}
\Lambda u= & \left.\left.\mathrm{i}(-1)^{p} \sum_{J, K, l} u_{J, K}\left(\zeta_{l}\right\lrcorner \zeta_{J}^{\star}\right) \wedge\left(\bar{\zeta}_{l}\right\lrcorner \bar{\zeta}_{K}^{\star}\right), \quad 1 \leqslant l \leqslant n, \\
\gamma \wedge u= & \mathrm{i}(-1)^{p} \sum_{J, K, m} \gamma_{m} u_{J, K} \zeta_{m}^{\star} \wedge \zeta_{J}^{\star} \wedge \bar{\zeta}_{m}^{\star} \wedge \bar{\zeta}_{K}^{\star}, \quad 1 \leqslant m \leqslant n, \\
{[\gamma, \Lambda] u=} & \left.\sum_{J, K, l, m} \gamma_{m} u_{J, K}\left(\left(\zeta_{l}^{\star} \wedge\left(\zeta_{m}\right\lrcorner \zeta_{J}^{\star}\right)\right) \wedge\left(\bar{\zeta}_{l}^{\star} \wedge\left(\bar{\zeta}_{m}\right\lrcorner \bar{\zeta}_{K}^{\star}\right)\right) \\
& \left.\left.\left.\quad-\left(\zeta_{m}\right\lrcorner\left(\zeta_{l}^{\star} \wedge \zeta_{J}^{\star}\right)\right) \wedge\left(\bar{\zeta}_{m}\right\lrcorner\left(\bar{\zeta}_{l}^{\star} \wedge \bar{\zeta}_{K}^{\star}\right)\right)\right) \\
= & \sum_{J, K, m} \gamma_{m} u_{J, K}\left(\zeta_{m}^{\star} \wedge\left(\zeta_{m}\right\lrcorner \zeta_{J}^{\star}\right) \wedge \bar{\zeta}_{K}^{\star} \\
& \left.\left.\quad+\zeta_{J}^{\star} \wedge \bar{\zeta}_{m}^{\star} \wedge\left(\bar{\zeta}_{m}\right\lrcorner \bar{\zeta}_{K}^{\star}\right)-\zeta_{J}^{\star} \wedge \bar{\zeta}_{K}^{\star}\right) \\
= & \sum_{J, K}\left(\sum_{m \in J} \gamma_{m}+\sum_{m \in K} \gamma_{m}-\sum_{1 \leqslant m \leqslant n} \gamma_{m}\right) u_{J, K} \zeta_{J}^{\star} \wedge \bar{\zeta}_{K}^{\star} .
\end{aligned}
$$

(5.9) Corollary. For every $u \in \Lambda^{p, q} T_{X}^{\star}$, we have

$$
[L, \Lambda] u=(p+q-n) u
$$

Proof. Indeed, if $\gamma=\omega$, we have $\gamma_{1}=\cdots=\gamma_{n}=1$.
This result can be generalized as follows: for every $u \in \Lambda^{k}\left(\mathbb{C} \otimes T_{X}\right)^{\star}$, we have

$$
\begin{equation*}
\left[L^{r}, \Lambda\right] u=r(k-n+r-1) L^{r-1} u \tag{5.10}
\end{equation*}
$$

In fact, it is clear that

$$
\begin{aligned}
{\left[L^{r}, \Lambda\right] u } & =\sum_{0 \leqslant m \leqslant r-1} L^{r-1-m}[L, \Lambda] L^{m} u \\
& =\sum_{0 \leqslant m \leqslant r-1}(2 m+k-n) L^{r-1-m} L^{m} u=(r(r-1)+r(k-n)) L^{r-1} u .
\end{aligned}
$$

## § 5.3. Primitive Elements and Hard Lefschetz Theorem

In this subsection, we prove a fundamental decomposition theorem for the representation of the unitary group $U\left(T_{X}\right) \simeq U(n)$ acting on the spaces $\Lambda^{p, q} T_{X}^{\star}$ of $(p, q)$-forms. It turns out that the representation is never irreducible if $0<p, q<n$.
(5.11) Definition. A homogeneous element $u \in \Lambda^{k}\left(\mathbb{C} \otimes T_{X}\right)^{\star}$ is called primitive if $\Lambda u=0$. The space of primitive elements of total degree $k$ will be denoted

$$
\operatorname{Prim}^{k} T_{X}^{\star}=\bigoplus_{p+q=k} \operatorname{Prim}^{p, q} T_{X}^{\star}
$$

Let $u \in \operatorname{Prim}^{k} T_{X}^{\star}$. Then

$$
\Lambda^{s} L^{r} u=\Lambda^{s-1}\left(\Lambda L^{r}-L^{r} \Lambda\right) u=r(n-k-r+1) \Lambda^{s-1} L^{r-1} u
$$

By induction, we get for $r \geqslant s$

$$
\begin{equation*}
\Lambda^{s} L^{r} u=r(r-1) \cdots(r-s+1) \cdot(n-k-r+1) \cdots(n-k-r+s) L^{r-s} u . \tag{5.12}
\end{equation*}
$$

Apply (5.12) for $r=n+1$. Then $L^{n+1} u$ is of degree $>2 n$ and therefore we have $L^{n+1} u=0$. This gives

$$
(n+1) \cdots(n+1-(s-1)) \cdot(-k)(-k+1) \cdots(-k+s-1) L^{n+1-s} u=0
$$

The integral coefficient is $\neq 0$ if $s \leqslant k$, hence:
(5.13) Corollary. If $u \in \operatorname{Prim}^{k} T_{X}^{\star}$, then $L^{s} u=0$ for $s \geqslant(n+1-k)_{+}$.
(5.14) Corollary. Prim $^{k} T_{X}^{\star}=0$ for $n+1 \leqslant k \leqslant 2 n$.

Proof. Apply Corollary 5.13 with $s=0$.
(5.15) Primitive decomposition formula. For every $u \in \Lambda^{k}\left(\mathbb{C} \otimes T_{X}\right)^{\star}$, there is a unique decomposition

$$
u=\sum_{r \geqslant(k-n)_{+}} L^{r} u_{r}, \quad u_{r} \in \operatorname{Prim}^{k-2 r} T_{X}^{\star}
$$

Furthermore $u_{r}=\Phi_{k, r}(L, \Lambda) u$ where $\Phi_{k, r}$ is a non commutative polynomial in $L, \Lambda$ with rational coefficients. As a consequence, there are direct sum decompositions of $U(n)$ representations

$$
\begin{aligned}
& \Lambda^{k}\left(\mathbb{C} \otimes T_{X}\right)^{\star}=\bigoplus_{r \geqslant(k-n)_{+}} L^{r} \operatorname{Prim}^{k-2 r} T_{X}^{\star}, \\
& \Lambda^{p, q} T_{X}^{\star}=\bigoplus_{r \geqslant(p+q-n)_{+}} L^{r} \operatorname{Prim}^{p-r, q-r} T_{X}^{\star}
\end{aligned}
$$

Proof of the uniqueness of the decomposition. Assume that $u=0$ and that $u_{r} \neq 0$ for some $r$. Let $s$ be the largest integer such that $u_{s} \neq 0$. Then

$$
\Lambda^{s} u=0=\sum_{(k-n)_{+} \leqslant r \leqslant s} \Lambda^{s} L^{r} u_{r}=\sum_{(k-n)_{+} \leqslant r \leqslant s} \Lambda^{s-r} \Lambda^{r} L^{r} u_{r} .
$$

But formula (5.12) shows that $\Lambda^{r} L^{r} u_{r}=c_{k, r} u_{r}$ for some non zero integral coefficient $c_{k, r}=r!(n-k+r+1) \cdots(n-k+2 r)$. Since $u_{r}$ is primitive we get $\Lambda^{s} L^{r} u_{r}=0$ when $r<s$, hence $u_{s}=0$, a contradiction.

Proof of the existence of the decomposition. We prove by induction on $s \geqslant(k-n)_{+}$that $\Lambda^{s} u=0$ implies

$$
\begin{equation*}
u=\sum_{(k-n)_{+} \leqslant r<s} L^{r} u_{r}, \quad u_{r}=\Phi_{k, r, s}(L, \Lambda) u \in \operatorname{Prim}^{k-2 r} T_{X}^{\star} . \tag{5.16}
\end{equation*}
$$

The Theorem will follow from the step $s=n+1$.
Assume that the result is true for $s$ and that $\Lambda^{s+1} u=0$. Then $\Lambda^{s} u$ is in $\operatorname{Prim}^{k-2 s} T_{X}^{\star}$. Since $s \geqslant(k-n)_{+}$we have $c_{k, s} \neq 0$ and we set

$$
\begin{aligned}
& u_{s}=\frac{1}{c_{k, s}} \Lambda^{s} u \in \operatorname{Prim}^{k-2 s} T_{X}^{\star} \\
& u^{\prime}=u-L^{s} u_{s}=\left(1-\frac{1}{c_{k, s}} L^{s} \Lambda^{s}\right) u
\end{aligned}
$$

By formula (5.12), we get

$$
\Lambda^{s} u^{\prime}=\Lambda^{s} u-\Lambda^{s} L^{s} u_{s}=\Lambda^{s} u-c_{k, s} u_{s}=0
$$

The induction hypothesis implies

$$
u^{\prime}=\sum_{(k-n)_{+} \leqslant r<s} L^{r} u_{r}^{\prime}, u_{r}^{\prime}=\Phi_{k, r, s}(L, \Lambda) u^{\prime} \in \operatorname{Prim}^{k-2 r} T_{X}^{\star}
$$

hence $u=\sum_{(k-n)_{+} \leqslant r \leqslant s} L^{r} u_{r}$ with

$$
\left\{\begin{array}{l}
u_{r}=u_{r}^{\prime}=\Phi_{k, r, s}(L, \Lambda)\left(1-\frac{1}{c_{k, s}} L^{s} \Lambda^{s}\right) u, r<s \\
u_{s}=\frac{1}{c_{k, s}} \Lambda^{s} u
\end{array}\right.
$$

It remains to prove the validity of the decomposition 5.16) for the initial step $s=(k-n)_{+}$, i.e. that $\Lambda^{s} u=0$ implies $u=0$. If $k \leqslant n$, then $s=0$ and there is nothing to prove. We are left with the case $k>n, \Lambda^{k-n} u=0$. Then $v=\star u \in \Lambda^{2 n-k}\left(\mathbb{C} \otimes T_{X}\right)^{\star}$ and $2 n-k<n$. Since the decomposition exists in degree $\leqslant n$ by what we have just proved, we get

$$
\begin{aligned}
& v=\star u=\sum_{r \geqslant 0} L^{r} v_{r}, v_{r} \in \operatorname{Prim}^{2 n-k-2 r} T_{X}^{\star}, \\
& 0=\star \Lambda^{k-n} u=L^{k-n} \star u=\sum_{r \geqslant 0} L^{r+k-n} v_{r},
\end{aligned}
$$

with degree $\left(L^{r+k-n} v_{r}\right)=2 n-k+2(k-n)=k$. The uniqueness part shows that $v_{r}=0$ for all $r$, hence $u=0$. The Theorem is proved.
(5.17) Corollary. The linear operators

$$
\begin{aligned}
L^{n-k}: \Lambda^{k}\left(\mathbb{C} \otimes T_{X}\right)^{\star} & \longrightarrow \Lambda^{2 n-k}\left(\mathbb{C} \otimes T_{X}\right)^{\star}, \\
L^{n-p-q}: \Lambda^{p, q} T_{X}^{\star} & \longrightarrow \Lambda^{n-q, n-p} T_{X}^{\star},
\end{aligned}
$$

are isomorphisms for all integers $k \leqslant n, p+q \leqslant n$.
Proof. For every $u \in \Lambda_{\mathbb{C}}^{k} T_{X}^{\star}$, the primitive decomposition $u=\sum_{r \geqslant 0} L^{r} u_{r}$ is mapped bijectively onto that of $L^{n-k} u$ :

$$
L^{n-k} u=\sum_{r \geqslant 0} L^{r+n-k} u_{r} .
$$

## §6. Commutation Relations

## §6.1. Commutation Relations on a Kähler Manifold

Assume first that $X=\Omega \subset \mathbb{C}^{n}$ is an open subset and that $\omega$ is the standard Kähler metric

$$
\omega=\mathrm{i} \sum_{1 \leqslant j \leqslant n} d z_{j} \wedge d \bar{z}_{j} .
$$

For any form $u \in \mathscr{C}^{\infty}\left(\Omega, \Lambda^{p, q} T_{X}^{\star}\right)$ we have

$$
\begin{align*}
d^{\prime} u & =\sum_{I, J, k} \frac{\partial u_{I, J}}{\partial z_{k}} d z_{k} \wedge d z_{I} \wedge d \bar{z}_{J} \\
d^{\prime \prime} u & =\sum_{I, J, k} \frac{\partial u_{I, J}}{\partial \bar{z}_{k}} d \bar{z}_{k} \wedge d z_{I} \wedge d \bar{z}_{J}
\end{align*}
$$

Since the global $L^{2}$ inner product is given by

$$
\langle\langle u, v\rangle\rangle=\int_{\Omega} \sum_{I, J} u_{I, J} \bar{v}_{I, J} d V,
$$

easy computations analogous to those of Example 3.12 show that

$$
\begin{align*}
d^{\prime \star} u & \left.=-\sum_{I, J, k} \frac{\partial u_{I, J}}{\partial \bar{z}_{k}} \frac{\partial}{\partial z_{k}}\right\lrcorner\left(d z_{I} \wedge d \bar{z}_{J}\right), \\
d^{\prime \prime \star} u & \left.=-\sum_{I, J, k} \frac{\partial u_{I, J}}{\partial z_{k}} \frac{\partial}{\partial \bar{z}_{k}}\right\lrcorner\left(d z_{I} \wedge d \bar{z}_{J}\right) .
\end{align*}
$$

We first prove a lemma due to [Akizuki and Nakano 1954].
(6.3) Lemma. In $\mathbb{C}^{n}$, we have $\left[d^{\prime \prime \star}, L\right]=\mathrm{i} d^{\prime}$.

Proof. Formula ( $6.2^{\prime \prime}$ ) can be written more briefly

$$
\left.d^{\prime \prime \star} u=-\sum_{k} \frac{\partial}{\partial \bar{z}_{k}}\right\lrcorner\left(\frac{\partial u}{\partial z_{k}}\right) .
$$

Then we get

$$
\left.\left.\left[d^{\prime \prime \star}, L\right] u=-\sum_{k} \frac{\partial}{\partial \bar{z}_{k}}\right\lrcorner\left(\frac{\partial}{\partial z_{k}}(\omega \wedge u)\right)+\omega \wedge \sum_{k} \frac{\partial}{\partial \bar{z}_{k}}\right\lrcorner\left(\frac{\partial u}{\partial z_{k}}\right) .
$$

Since $\omega$ has constant coefficients, we have $\frac{\partial}{\partial z_{k}}(\omega \wedge u)=\omega \wedge \frac{\partial u}{\partial z_{k}}$ and therefore

$$
\begin{aligned}
{\left[d^{\prime \prime \star}, L\right] u } & \left.\left.=-\sum_{k}\left(\frac{\partial}{\partial \bar{z}_{k}}\right\lrcorner\left(\omega \wedge \frac{\partial u}{\partial z_{k}}\right)-\omega \wedge\left(\frac{\partial}{\partial \bar{z}_{k}}\right\lrcorner \frac{\partial u}{\partial z_{k}}\right)\right) \\
& \left.=-\sum_{k}\left(\frac{\partial}{\partial \bar{z}_{k}}\right\lrcorner \omega\right) \wedge \frac{\partial u}{\partial z_{k}} .
\end{aligned}
$$

Clearly $\left.\frac{\partial}{\partial \bar{z}_{k}}\right\lrcorner \omega=-\mathrm{i} d z_{k}$, so

$$
\left[d^{\prime \prime \star}, L\right] u=\mathrm{i} \sum_{k} d z_{k} \wedge \frac{\partial u}{\partial z_{k}}=\mathrm{i} d^{\prime} u
$$

We are now ready to derive the basic commutation relations in the case of an arbitrary Kähler manifold $(X, \omega)$.
(6.4) Theorem. If $(X, \omega)$ is Kähler, then

$$
\begin{array}{ll}
{\left[d^{\prime \prime \star}, L\right]=\mathrm{i} d^{\prime},} & {\left[d^{\prime \star}, L\right]=-\mathrm{i} d^{\prime \prime},} \\
{\left[\Lambda, d^{\prime \prime}\right]=-\mathrm{i} d^{\prime \star},} & \\
\left.\hline \Lambda, d^{\prime}\right]=\mathrm{i} d^{\prime \prime \star} .
\end{array}
$$

Proof. It is sufficient to verify the first relation, because the second one is the conjugate of the first, and the relations of the second line are the adjoint of those of the first line. If $\left(z_{j}\right)$ is a geodesic coordinate system at a point $x_{0} \in X$, then for any $(p, q)$-forms $u, v$ with compact support in a neighborhood of $x_{0}$, (4.9) implies

$$
\langle\langle u, v\rangle\rangle=\int_{M}\left(\sum_{I, J} u_{I J} \bar{v}_{I J}+\sum_{I, J, K, L} a_{I J K L} u_{I J} \bar{v}_{K L}\right) d V,
$$

with $a_{I J K L}(z)=O\left(|z|^{2}\right)$ at $x_{0}$. An integration by parts as in (3.12) and (6.2") yields

$$
\left.d^{\prime \prime \star} u=-\sum_{I, J, k} \frac{\partial u_{I, J}}{\partial z_{k}} \frac{\partial}{\partial \bar{z}_{k}}\right\lrcorner\left(d z_{I} \wedge d \bar{z}_{J}\right)+\sum_{I, J, K, L} b_{I J K L} u_{I J} d z_{K} \wedge d \bar{z}_{L}
$$

where the coefficients $b_{I J K L}$ are obtained by derivation of the $a_{I J K L}$ 's. Therefore $b_{I J K L}=O(|z|)$. Since $\partial \omega / \partial z_{k}=O(|z|)$, the proof of Lemma 6.3 implies here $\left[d^{\prime \prime \star}, L\right] u=$ $\mathrm{id} d^{\prime} u+O(|z|)$, in particular both terms coincide at every given point $x_{0} \in X$.
(6.5) Corollary. If $(X, \omega)$ is Kähler, the complex Laplace-Beltrami operators satisfy

$$
\Delta^{\prime}=\Delta^{\prime \prime}=\frac{1}{2} \Delta .
$$

Proof. It will be first shown that $\Delta^{\prime \prime}=\Delta^{\prime}$. We have

$$
\Delta^{\prime \prime}=\left[d^{\prime \prime}, d^{\prime \prime \star}\right]=-\mathrm{i}\left[d^{\prime \prime},\left[\Lambda, d^{\prime}\right]\right] .
$$

Since $\left[d^{\prime}, d^{\prime \prime}\right]=0$, Jacobi's identity (5.7) implies

$$
-\left[d^{\prime \prime},\left[\Lambda, d^{\prime}\right]\right]+\left[d^{\prime},\left[d^{\prime \prime}, \Lambda\right]\right]=0
$$

hence $\Delta^{\prime \prime}=\left[d^{\prime},-\mathrm{i}\left[d^{\prime \prime}, \Lambda\right]\right]=\left[d^{\prime}, d^{\prime \star}\right]=\Delta^{\prime}$. On the other hand

$$
\Delta=\left[d^{\prime}+d^{\prime \prime}, d^{\prime \star}+d^{\prime \prime \star}\right]=\Delta^{\prime}+\Delta^{\prime \prime}+\left[d^{\prime}, d^{\prime \prime \star}\right]+\left[d^{\prime \prime}, d^{\prime \star}\right] .
$$

Thus, it is enough to prove:
(6.6) Lemma. $\left[d^{\prime}, d^{\prime \prime \star}\right]=0,\left[d^{\prime \prime}, d^{\prime \star}\right]=0$.

Proof. We have $\left[d^{\prime}, d^{\prime \prime *}\right]=-\mathrm{i}\left[d^{\prime},\left[\Lambda, d^{\prime}\right]\right]$ and (5.7) implies

$$
-\left[d^{\prime},\left[\Lambda, d^{\prime}\right]\right]+\left[\Lambda,\left[d^{\prime}, d^{\prime}\right]\right]+\left[d^{\prime},\left[d^{\prime}, \Lambda\right]\right]=0
$$

hence $-2\left[d^{\prime},\left[\Lambda, d^{\prime}\right]\right]=0$ and $\left[d^{\prime}, d^{\prime \prime \star}\right]=0$. The second relation $\left[d^{\prime \prime}, d^{\prime \star}\right]=0$ is the adjoint of the first.
(6.7) Theorem. $\Delta$ commutes with all operators $\star, d^{\prime}, d^{\prime \prime}, d^{\prime \star}, d^{\prime \prime \star}, L, \Lambda$.

Proof. The identities $\left[d^{\prime}, \Delta^{\prime}\right]=\left[d^{\prime \star}, \Delta^{\prime}\right]=0,\left[d^{\prime \prime}, \Delta^{\prime \prime}\right]=\left[d^{\prime \prime \star}, \Delta^{\prime \prime}\right]=0$ and $[\Delta, \star]=0$ are immediate. Furthermore, the equality $\left[d^{\prime}, L\right]=d^{\prime} \omega=0$ together with the Jacobi identity implies

$$
\left[L, \Delta^{\prime}\right]=\left[L,\left[d^{\prime}, d^{\prime \star}\right]\right]=-\left[d^{\prime},\left[d^{\prime \star}, L\right]\right]=\mathrm{i}\left[d^{\prime}, d^{\prime \prime}\right]=0 .
$$

By adjunction, we also get $\left[\Delta^{\prime}, \Lambda\right]=0$.

## § 6.2. Commutation Relations on Hermitian Manifolds

We are going to extend the commutation relations of $\S 6.1$ to an arbitrary hermitian manifold $(X, \omega)$. In that case $\omega$ is no longer tangent to a constant metric, and the commutation relations involve extra terms arising from the torsion of $\omega$. Theorem 6.8 below is taken from [Demailly 1984], but the idea was already contained in [Griffiths 1966].
(6.8) Theorem. Let $\tau$ be the operator of type $(1,0)$ and order 0 defined by $\tau=\left[\Lambda, d^{\prime} \omega\right]$. Then
a) $\quad\left[d^{\prime \prime \star}, L\right]=\mathrm{i}\left(d^{\prime}+\tau\right)$,
b) $\quad\left[d^{\prime \star}, L\right]=-\mathrm{i}\left(d^{\prime \prime}+\bar{\tau}\right)$,
c) $\left[\Lambda, d^{\prime \prime}\right]=-\mathrm{i}\left(d^{\prime \star}+\tau^{\star}\right)$,
d) $\quad\left[\Lambda, d^{\prime}\right]=\mathrm{i}\left(d^{\prime \prime \star}+\bar{\tau}^{\star}\right)$;
$d^{\prime} \omega$ will be called the torsion form of $\omega$, and $\tau$ the torsion operator.
Proof. b) follows from a) by conjugation, whereas c), d) follow from a), b) by adjunction. It is therefore enough to prove relation a).

Let $\left(z_{j}\right)_{1 \leqslant j \leqslant n}$ be complex coordinates centered at a point $x_{0} \in X$, such that the $n$ tuple $\left(\partial / \partial z_{1}, \ldots, \partial / \partial z_{n}\right)$ is an orthonormal basis of $T_{x_{0}} X$ for the metric $\omega\left(x_{0}\right)$. Consider the metric with constant coefficients

$$
\omega_{0}=\mathrm{i} \sum_{1 \leqslant j \leqslant n} d z_{j} \wedge d \bar{z}_{j} .
$$

The metric $\omega$ can then be written

$$
\omega=\omega_{0}+\gamma \text { with } \gamma=O(|z|) .
$$

Denote by $\langle,\rangle_{0}, L_{0}, \Lambda_{0}, d_{0}^{\prime \star}, d_{0}^{\prime \prime \star}$ the inner product and the operators associated to the constant metric $\omega_{0}$, and let $d V_{0}=\omega_{0}^{n} / 2^{n} n$ !. The proof of relation a) is based on a Taylor expansion of $L, \Lambda, d^{\prime \star}, d^{\prime \prime \star}$ in terms of the operators with constant coefficients $L_{0}, \Lambda_{0}, d_{0}^{\prime \star}, d_{0}^{\prime \prime *}$.
(6.9) Lemma. Let $u, v \in \mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star}\right)$. Then in a neighborhood of $x_{0}$

$$
\langle u, v\rangle d V=\left\langle u-\left[\gamma, \Lambda_{0}\right] u, v\right\rangle_{0} d V_{0}+O\left(|z|^{2}\right)
$$

Proof. In a neighborhood of $x_{0}$, let

$$
\gamma=\mathrm{i} \sum_{1 \leqslant j \leqslant n} \gamma_{j} \zeta_{j}^{\star} \wedge \bar{\zeta}_{j}^{\star}, \quad \gamma_{1} \leqslant \gamma_{2} \leqslant \cdots \leqslant \gamma_{n}
$$

be a diagonalization of the $(1,1)$-form $\gamma(z)$ with respect to an orthonormal basis $\left(\zeta_{j}\right)_{1 \leqslant j \leqslant n}$ of $T_{z} X$ for $\omega_{0}(z)$. We thus have

$$
\omega=\omega_{0}+\gamma=\mathrm{i} \sum \lambda_{j} \zeta_{j}^{\star} \wedge \bar{\zeta}_{j}^{\star}
$$

with $\lambda_{j}=1+\gamma_{j}$ and $\gamma_{j}=O(|z|)$. Set now

$$
\begin{gathered}
J=\left\{j_{1}, \ldots, j_{p}\right\}, \quad \zeta_{J}^{\star}=\zeta_{j_{1}}^{\star} \wedge \cdots \wedge \zeta_{j_{p}}^{\star}, \quad \lambda_{J}=\lambda_{j_{1}} \cdots \lambda_{j_{p}}, \\
u=\sum u_{J, K} \zeta_{J}^{\star} \wedge \bar{\zeta}_{K}^{\star}, \quad v=\sum v_{J, K} \zeta_{J}^{\star} \wedge \bar{\zeta}_{K}^{\star}
\end{gathered}
$$

where summations are extended to increasing multi-indices $J, K$ such that $|J|=p$, $|K|=q$. With respect to $\omega$ we have $\left\langle\zeta_{j}^{\star}, \zeta_{j}^{\star}\right\rangle=\lambda_{j}^{-1}$, hence

$$
\begin{aligned}
\langle u, v\rangle d V & =\sum_{J, K} \lambda_{J}^{-1} \lambda_{K}^{-1} u_{J, K} \bar{v}_{J, K} \lambda_{1} \cdots \lambda_{n} d V_{0} \\
& =\sum_{J, K}\left(1-\sum_{j \in J} \gamma_{j}-\sum_{j \in K} \gamma_{j}+\sum_{1 \leqslant j \leqslant n} \gamma_{j}\right) u_{J, K} \bar{v}_{J, K} d V_{0}+O\left(|z|^{2}\right) .
\end{aligned}
$$

Lemma 6.9 follows if we take Prop. 5.8 into account.
(6.10) Lemma. $d^{\prime \prime \star}=d_{0}^{\prime \prime \star}+\left[\Lambda_{0},\left[d_{0}^{\prime \prime \star}, \gamma\right]\right]$ at point $x_{0}$, i.e. at this point both operators have the same formal expansion.

Proof. Since $d^{\prime \prime *}$ is an operator of order 1, Lemma 6.9 shows that $d^{\prime \prime *}$ coincides at point $x_{0}$ with the formal adjoint of $d^{\prime \prime}$ for the metric

$$
\langle\langle u, v\rangle\rangle_{1}=\int_{X}\left\langle u-\left[\gamma, \Lambda_{0}\right] u, v\right\rangle_{0} d V_{0} .
$$

For any compactly supported $u \in \mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star}\right), v \in \mathscr{C}^{\infty}\left(X, \Lambda^{p, q-1} T_{X}^{\star}\right)$ we get by definition

$$
\left\langle\left\langle u, d^{\prime \prime} v\right\rangle_{1}=\int_{X}\left\langle u-\left[\gamma, \Lambda_{0}\right] u, d^{\prime \prime} v\right\rangle_{0} d V_{0}=\int_{X}\left\langle d_{0}^{\prime \prime \star} u-d_{0}^{\prime \prime \star}\left[\gamma, \Lambda_{0}\right] u, v\right\rangle_{0} d V_{0} .\right.
$$

Since $\omega$ and $\omega_{0}$ coincide at point $x_{0}$ and since $\gamma\left(x_{0}\right)=0$ we obtain at this point

$$
\begin{aligned}
d^{\prime \prime \star} u & =d_{0}^{\prime \prime \star} u-d_{0}^{\prime \prime \star}\left[\gamma, \Lambda_{0}\right] u=d_{0}^{\prime \prime \star} u-\left[d_{0}^{\prime \prime \star},\left[\gamma, \Lambda_{0}\right]\right] u \\
d^{\prime \prime \star} & =d_{0}^{\prime \prime \star}-\left[d_{0}^{\prime \prime \star},\left[\gamma, \Lambda_{0}\right]\right] .
\end{aligned}
$$

We have $\left[\Lambda_{0}, d_{0}^{\prime \prime \star}\right]=\left[d^{\prime \prime}, L_{0}\right]^{\star}=0$ since $d^{\prime \prime} \omega_{0}=0$. The Jacobi identity (5.7) implies

$$
\left[d_{0}^{\prime \prime \star},\left[\gamma, \Lambda_{0}\right]\right]+\left[\Lambda_{0},\left[d_{0}^{\prime \prime \star}, \gamma\right]\right]=0
$$

and Lemma 6.10 follows.
Proof Proof of formula 6.8 a). The equality $L=L_{0}+\gamma$ and Lemma 6.10 yield

$$
\begin{equation*}
\left[L, d^{\prime \prime \star}\right]=\left[L_{0}, d_{0}^{\prime \prime \star}\right]+\left[L_{0},\left[\Lambda_{0},\left[d_{0}^{\prime \prime \star}, \gamma\right]\right]\right]+\left[\gamma, d_{0}^{\prime \prime \star}\right] \tag{6.11}
\end{equation*}
$$

at point $x_{0}$, because the triple bracket involving $\gamma$ twice vanishes at $x_{0}$. From the Jacobi identity applied to $C=\left[d_{0}^{\prime \prime \star}, \gamma\right]$, we get

$$
\left\{\begin{array}{l}
{\left[L_{0},\left[\Lambda_{0}, C\right]\right]=-\left[\Lambda_{0},\left[C, L_{0}\right]\right]-\left[C,\left[L_{0}, \Lambda_{0}\right]\right]}  \tag{6.12}\\
{\left[C, L_{0}\right]=\left[L_{0},\left[d_{0}^{\prime \prime \star}, \gamma\right]\right]=\left[\gamma,\left[L_{0}, d_{0}^{\prime \prime \star}\right]\right] \quad\left(\text { since }\left[\gamma, L_{0}\right]=0\right)}
\end{array}\right.
$$

Lemma 6.3 yields $\left[L_{0}, d_{0}^{\prime \prime *}\right]=-\mathrm{i} d^{\prime}$, hence

$$
\begin{equation*}
\left[C, L_{0}\right]=-\left[\gamma, \mathrm{i} d^{\prime}\right]=\mathrm{i} d^{\prime} \gamma=\mathrm{i} d^{\prime} \omega \tag{6.13}
\end{equation*}
$$

On the other hand, $C$ is of type $(1,0)$ and Cor. 5.9 gives

$$
\begin{equation*}
\left[C,\left[L_{0}, \Lambda_{0}\right]\right]=-C=-\left[d_{0}^{\prime \prime \star}, \gamma\right] . \tag{6.14}
\end{equation*}
$$

From (6.12), (6.13), (6.14) we get

$$
\left[L_{0},\left[\Lambda_{0},\left[d_{0}^{\prime \prime \star}, \gamma\right]\right]\right]=-\left[\Lambda_{0}, \mathrm{id} d^{\prime} \omega\right]+\left[d_{0}^{\prime \prime \star}, \gamma\right]
$$

This last equality combined with (6.11) implies

$$
\left[L, d^{\prime \prime *}\right]=\left[L_{0}, d_{0}^{\prime \prime *}\right]-\left[\Lambda_{0}, \mathrm{i} d^{\prime} \omega\right]=-\mathrm{i}\left(d^{\prime}+\tau\right)
$$

at point $x_{0}$. Formula 6.8 a) is proved.
(6.15) Corollary. The complex Laplace-Beltrami operators satisfy

$$
\begin{aligned}
\Delta^{\prime \prime} & =\Delta^{\prime}+\left[d^{\prime}, \tau^{\star}\right]-\left[d^{\prime \prime}, \bar{\tau}^{\star}\right], \\
{\left[d^{\prime}, d^{\prime \prime \star}\right] } & =-\left[d^{\prime}, \bar{\tau}^{\star}\right], \quad\left[d^{\prime \prime}, d^{\prime \star}\right]=-\left[d^{\prime \prime}, \tau^{\star}\right], \\
\Delta & =\Delta^{\prime}+\Delta^{\prime \prime}-\left[d^{\prime}, \bar{\tau}^{\star}\right]-\left[d^{\prime \prime}, \tau^{\star}\right] .
\end{aligned}
$$

Therefore $\Delta^{\prime}, \Delta^{\prime \prime}$ and $\frac{1}{2} \Delta$ no longer coincide, but they differ by linear differential operators of order 1 only.

Proof. As in the Kähler case (Cor. 6.5 and Lemma 6.6), we find

$$
\begin{aligned}
\Delta^{\prime \prime} & =\left[d^{\prime \prime}, d^{\prime \prime \star}\right]=\left[d^{\prime \prime},-i\left[\Lambda, d^{\prime}\right]-\bar{\tau}^{\star}\right] \\
& =\left[d^{\prime},-\mathrm{i}\left[d^{\prime \prime}, \Lambda\right]\right]-\left[d^{\prime \prime}, \bar{\tau}^{\star}\right]=\Delta^{\prime}+\left[d^{\prime}, \tau^{\star}\right]-\left[d^{\prime \prime}, \bar{\tau}^{\star}\right], \\
{\left[d^{\prime}, d^{\prime \prime \star}+\bar{\tau}^{\star}\right] } & =-\mathrm{i}\left[d^{\prime},\left[\Lambda, d^{\prime}\right]\right]=0,
\end{aligned}
$$

and the first two lines are proved. The third one is an immediate consequence of the second.

## $\S$ 7. Groups $\mathscr{H}^{p, q}(X, E)$ and Serre Duality

Let $(X, \omega)$ be a compact hermitian manifold and $E$ a holomorphic hermitian vector bundle of rank $r$ over $X$. We denote by $D_{E}$ the Chern connection of $E$, by $D_{E}^{\star}=-\star D_{E} \star$ the formal adjoint of $D_{E}$, and by $D_{E}^{\prime \star}, D_{E}^{\prime \prime \star}$ the components of $D_{E}^{\star}$ of type $(-1,0)$ and $(0,-1)$.

Corollary 6.8 implies that the principal part of the operator $\Delta_{E}^{\prime \prime}=D^{\prime \prime} D_{E}^{\prime \prime \star}+D_{E}^{\prime \prime \star} D^{\prime \prime}$ is one half that of $\Delta_{E}$. The operator $\Delta_{E}^{\prime \prime}$ acting on each space $\mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star} \otimes E\right)$ is thus a self-adjoint elliptic operator. Since $D^{\prime \prime 2}=0$, the following results can be obtained in a way similar to those of $\S 3.3$.
(7.1) Theorem. For every bidegree ( $p, q$ ), there exists an orthogonal decomposition

$$
\mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star} \otimes E\right)=\mathscr{H}^{p, q}(X, E) \oplus \operatorname{Im} D_{E}^{\prime \prime} \oplus \operatorname{Im} D_{E}^{\prime \prime \star}
$$

where $\mathscr{H}^{p, q}(X, E)$ is the space of $\Delta_{E}^{\prime \prime}$-harmonic forms in $\mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star} \otimes E\right)$.
The above decomposition shows that the subspace of $d^{\prime \prime}$-cocycles in $\mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star} \otimes\right.$ $E)$ is $\mathscr{H}^{p, q}(X, E) \oplus \operatorname{Im} D_{E}^{\prime \prime}$. From this, we infer
(7.2) Hodge isomorphism theorem. The Dolbeault cohomology group $H^{p, q}(X, E)$ is finite dimensional, and there is an isomorphism

$$
H^{p, q}(X, E) \simeq \mathscr{H}^{p, q}(X, E)
$$

(7.3) Serre duality theorem. The bilinear pairing

$$
H^{p, q}(X, E) \times H^{n-p, n-q}\left(X, E^{\star}\right) \longrightarrow \mathbb{C}, \quad(s, t) \longmapsto \int_{M} s \wedge t
$$

is a non degenerate duality.
Proof. Let $s_{1} \in \mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star} \otimes E\right), s_{2} \in \mathscr{C}^{\infty}\left(X, \Lambda^{n-p, n-q-1} T_{X}^{\star} \otimes E\right)$. Since $s_{1} \wedge s_{2}$ is of bidegree ( $n, n-1$ ), we have

$$
\begin{equation*}
d\left(s_{1} \wedge s_{2}\right)=d^{\prime \prime}\left(s_{1} \wedge s_{2}\right)=d^{\prime \prime} s_{1} \wedge s_{2}+(-1)^{p+q} s_{1} \wedge d^{\prime \prime} s_{2} \tag{7.4}
\end{equation*}
$$

Stokes' formula implies that the above bilinear pairing can be factorized through Dolbeault cohomology groups. The \# operator defined in $\S 3.1$ is such that

$$
\#: \mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star} \otimes E\right) \longrightarrow \mathscr{C}^{\infty}\left(X, \Lambda^{n-p, n-q} T_{X}^{\star} \otimes E^{\star}\right)
$$

Furthermore, (3.20) implies

$$
\begin{gathered}
d^{\prime \prime}(\# s)=(-1)^{\operatorname{deg} s} \# D_{E}^{\prime \prime \star} s, \quad D_{E^{\star}}^{\prime \prime \star}(\# s)=(-1)^{\operatorname{deg} s+1} \# D_{E}^{\prime \prime \star} s \\
\Delta_{E^{\star}}^{\prime \prime}(\# s)=\# \Delta_{E}^{\prime \prime} s
\end{gathered}
$$

where $D_{E^{\star}}$ is the Chern connection of $E^{\star}$. Consequently, $s \in \mathscr{H}^{p, q}(X, E)$ if and only if $\# s \in \mathscr{H}^{n-p, n-q}\left(X, E^{\star}\right)$. Theorem 7.3 is then a consequence of the fact that the integral $\|s\|^{2}=\int_{X} s \wedge \# s$ does not vanish unless $s=0$.

## § 8. Cohomology of Compact Kähler Manifolds

## $\S$ 8.1. Bott-Chern Cohomology Groups

Let $X$ be for the moment an arbitrary complex manifold. The following "cohomology" groups are helpful to describe Hodge theory on compact complex manifolds which are not necessarily Kähler.
(8.1) Definition. We define the Bott-Chern cohomology groups of $X$ to be

$$
H_{\mathrm{BC}}^{p, q}(X, \mathbb{C})=\left(\mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star}\right) \cap \operatorname{ker} d\right) / d^{\prime} d^{\prime \prime} \mathscr{C}^{\infty}\left(X, \Lambda^{p-1, q-1} T_{X}^{\star}\right) .
$$

Then $H_{\mathrm{BC}}^{\bullet \bullet \bullet}(X, \mathbb{C})$ has the structure of a bigraded algebra, which we call the Bott-Chern cohomology algebra of $X$.

As the group $d^{\prime} d^{\prime \prime} \mathscr{C}^{\infty}\left(X, \Lambda^{p-1, q-1} T_{X}^{\star}\right)$ is contained in both coboundary groups $d^{\prime \prime} \mathscr{C}^{\infty}\left(X, \Lambda^{p, q-1} T_{X}^{\star}\right)$ or $d^{\infty} \mathscr{C}^{\infty}\left(X, \Lambda^{p+q-1}\left(\mathbb{C} \otimes T_{X}\right)^{\star}\right)$, there are canonical morphisms

$$
\begin{align*}
& H_{\mathrm{BC}}^{p, q}(X, \mathbb{C}) \longrightarrow H^{p, q}(X, \mathbb{C})  \tag{8.2}\\
& H_{\mathrm{BC}}^{p, q}(X, \mathbb{C}) \longrightarrow H_{\mathrm{DR}}^{p+q}(X, \mathbb{C}), \tag{8.3}
\end{align*}
$$

of the Bott-Chern cohomology to the Dolbeault or De Rham cohomology. These morphisms are homomorphisms of $\mathbb{C}$-algebras. It is also clear from the definition that we have the symmetry property $H_{\mathrm{BC}}^{q, p}(X, \mathbb{C})=\overline{H_{\mathrm{BC}}^{p, q}(X, \mathbb{C})}$. It can be shown from the HodgeFrölicher spectral sequence (see § 11 and Exercise 13.??) that $H_{\mathrm{BC}}^{p, q}(X, \mathbb{C})$ is always finite dimensional if $X$ is compact.

## § 8.2. Hodge Decomposition Theorem

We suppose from now on that $(X, \omega)$ is a compact Kähler manifold. The equality $\Delta=2 \Delta^{\prime \prime}$ shows that $\Delta$ is homogeneous with respect to bidegree and that there is an orthogonal decomposition

$$
\begin{equation*}
\mathscr{H}^{k}(X, \mathbb{C})=\bigoplus_{p+q=k} \mathscr{H}^{p, q}(X, \mathbb{C}) \tag{8.4}
\end{equation*}
$$

As $\overline{\Delta^{\prime \prime}}=\Delta^{\prime}=\Delta^{\prime \prime}$, we also have $\mathscr{H}^{q, p}(X, \mathbb{C})=\overline{\mathscr{H}^{p, q}(X, \mathbb{C})}$. Using the Hodge isomorphism theorems for the De Rham and Dolbeault cohomology, we get:
(8.5) Hodge decomposition theorem. On a compact Kähler manifold, there are canonical isomorphisms

$$
\begin{array}{ll}
H^{k}(X, \mathbb{C}) \simeq \bigoplus_{p+q=k} H^{p, q}(X, \mathbb{C}) & (\text { Hodge decomposition }), \\
H^{q, p}(X, \mathbb{C}) \simeq \overline{H^{p, q}}(X, \mathbb{C}) & (\text { Hodge symmetry })
\end{array}
$$

The only point which is not a priori completely clear is that this decomposition is independent of the Kähler metric. In order to show that this is the case, one can use the following Lemma, which allows us to compare all three types of cohomology groups considered in § 8.1.
(8.6) Lemma. Let $u$ be a d-closed $(p, q)$-form. The following properties are equivalent:
a) $u$ is d-exact;
$\left.\mathrm{b}^{\prime}\right) \mathrm{u}$ is $d^{\prime}$-exact;
$\left.\mathrm{b}^{\prime \prime}\right) u$ is $d^{\prime \prime}$-exact;
c) $u$ is $d^{\prime} d^{\prime \prime}$-exact, i.e. $u$ can be written $u=d^{\prime} d^{\prime \prime} v$.
d) $u$ is orthogonal to $\mathscr{H}^{p, q}(X, \mathbb{C})$.

Proof. It is obvious that $c$ ) implies a ), $\left.\mathrm{b}^{\prime}\right), \mathrm{b}^{\prime \prime}$ ) and that a) or $\mathrm{b}^{\prime}$ ) or $\mathrm{b}^{\prime \prime}$ ) implies $d$ ). It is thus sufficient to prove that d) implies c). As $d u=0$, we have $d^{\prime} u=d^{\prime \prime} u=0$, and as $u$ is supposed to be orthogonal to $\mathscr{H}^{p, q}(X, \mathbb{C})$, Th. 7.1 implies $u=d^{\prime \prime} s, s \in \mathscr{C}^{\infty}\left(X, \Lambda^{p, q-1} T_{X}^{\star}\right)$. By the analogue of Th. 7.1 for $d^{\prime}$, we have $s=h+d^{\prime} v+d^{\prime \star} w$, with $h \in \mathscr{H}^{p, q-1}(X, \mathbb{C})$, $v \in \mathscr{C}^{\infty}\left(X, \Lambda^{p-1, q-1} T_{X}^{\star}\right)$ and $w \in \mathscr{C}^{\infty}\left(X, \Lambda^{p+1, q-1} T_{X}^{\star}\right)$. Therefore

$$
u=d^{\prime \prime} d^{\prime} v+d^{\prime \prime} d^{\prime \star} w=-d^{\prime} d^{\prime \prime} v-d^{\prime \star} d^{\prime \prime} w
$$

in view of Lemma 6.6. As $d^{\prime} u=0$, the component $d^{\prime \star} d^{\prime \prime} w$ orthogonal to ker $d^{\prime}$ must be zero.

From Lemma 8.6 we infer the following Corollary, which in turn implies that the Hodge decomposition does not depend on the Kähler metric.
(8.7) Corollary. Let $X$ be a compact Kähler manifold. Then the natural morphisms

$$
H_{\mathrm{BC}}^{p, q}(X, \mathbb{C}) \longrightarrow H^{p, q}(X, \mathbb{C}), \quad \bigoplus_{p+q=k} H_{\mathrm{BC}}^{p, q}(X, \mathbb{C}) \longrightarrow H_{\mathrm{DR}}^{k}(X, \mathbb{C})
$$

are isomorphisms.
Proof. The surjectivity of $H_{\mathrm{BC}}^{p, q}(X, \mathbb{C}) \rightarrow H^{p, q}(X, \mathbb{C})$ comes from the fact that every class in $H^{p, q}(X, \mathbb{C})$ can be represented by a harmonic $(p, q)$-form, thus by a $d$-closed $(p, q)$ form; the injectivity means nothing more than the equivalence $\left(8.5 \mathrm{~b}^{\prime \prime}\right) \Leftrightarrow(8.5 \mathrm{c})$. Hence $H_{\mathrm{BC}}^{p, q}(X, \mathbb{C}) \simeq H^{p, q}(X, \mathbb{C}) \simeq \mathscr{H}^{p, q}(X, \mathbb{C})$, and the isomorphism $\bigoplus_{p+q=k} H_{\mathrm{BC}}^{p, q}(X, \mathbb{C}) \longrightarrow$ $H_{\mathrm{DR}}^{k}(X, \mathbb{C})$ follows from (8.4).

Let us quote now two simple applications of Hodge theory. The first of these is a computation of the Dolbeault cohomology groups of $\mathbb{P}^{n}$. As $H_{\mathrm{DR}}^{2 p}\left(\mathbb{P}^{n}, \mathbb{C}\right)=\mathbb{C}$ and $H^{p, p}\left(\mathbb{P}^{n}, \mathbb{C}\right) \ni\left\{\omega^{p}\right\} \neq 0$, the Hodge decomposition formula implies:
(8.8) Application. The Dolbeault cohomology groups of $\mathbb{P}^{n}$ are

$$
H^{p, p}\left(\mathbb{P}^{n}, \mathbb{C}\right)=\mathbb{C} \quad \text { for } 0 \leqslant p \leqslant n, \quad H^{p, q}\left(\mathbb{P}^{n}, \mathbb{C}\right)=0 \quad \text { for } p \neq q .
$$

(8.9) Proposition. Every holomorphic p-form on a compact Kähler manifold $X$ is d-closed.

Proof. If $u$ is a holomorphic form of type $(p, 0)$ then $d^{\prime \prime} u=0$. Furthermore $d^{\prime \prime \star} u$ is of type ( $p,-1$ ), hence $d^{\prime \prime \star} u=0$. Therefore $\Delta u=2 \Delta^{\prime \prime} u=0$, which implies $d u=0$.
(8.10) Example. Consider the Heisenberg group $G \subset \mathrm{Gl}_{3}(\mathbb{C})$, defined as the subgroup of matrices

$$
M=\left(\begin{array}{lll}
1 & x & z \\
0 & 1 & y \\
0 & 0 & 1
\end{array}\right) \quad, \quad(x, y, z) \in \mathbb{C}^{3}
$$

Let $\Gamma$ be the discrete subgroup of matrices with entries $x, y, z \in \mathbb{Z}[\mathrm{i}]$ (or more generally in the ring of integers of an imaginary quadratic field). Then $X=G / \Gamma$ is a compact complex 3 -fold, known as the Iwasawa manifold. The equality

$$
M^{-1} d M=\left(\begin{array}{ccc}
0 & d x & d z-x d y \\
0 & 0 & d y \\
0 & 0 & 0
\end{array}\right)
$$

shows that $d x, d y, d z-x d y$ are left invariant holomorphic 1-forms on $G$. These forms induce holomorphic 1-forms on the quotient $X=G / \Gamma$. Since $d z-x d y$ is not $d$-closed, we see that $X$ cannot be Kähler.

## § 8.3. Primitive Decomposition and Hard Lefschetz Theorem

We first introduce some standard notation. The Betti numbers and Hodge numbers of $X$ are by definition

$$
\begin{equation*}
b_{k}=\operatorname{dim}_{\mathbb{C}} H^{k}(X, \mathbb{C}), \quad h^{p, q}=\operatorname{dim}_{\mathbb{C}} H^{p, q}(X, \mathbb{C}) \tag{8.11}
\end{equation*}
$$

Thanks to Hodge decomposition, these numbers satisfy the relations

$$
\begin{equation*}
b_{k}=\sum_{p+q=k} h^{p, q}, \quad h^{q, p}=h^{p, q} . \tag{8.12}
\end{equation*}
$$

As a consequence, the Betti numbers $b_{2 k+1}$ of a compact Kähler manifold are even. Note that the Serre duality theorem gives the additional relation $h^{p, q}=h^{n-p, n-q}$, which holds as soon as $X$ is compact. The existence of primitive decomposition implies other interesting specific features of the cohomology algebra of compact Kähler manifolds.
(8.13) Lemma. If $u=\sum_{r \geqslant(k-n)_{+}} L^{r} u_{r}$ is the primitive decomposition of a harmonic $k$-form $u$, then all components $u_{r}$ are harmonic.

Proof. Since $[\Delta, L]=0$, we get $0=\Delta u=\sum_{r} L^{r} \Delta u_{r}$, hence $\Delta u_{r}=0$ by uniqueness.
Let us denote by $\operatorname{Prim} \mathscr{H}^{k}(X, \mathbb{C})=\bigoplus_{p+q=k} \operatorname{Prim} \mathscr{H}^{p, q}(X, \mathbb{C})$ the spaces of primitive harmonic $k$-forms and let $b_{k \text {,prim }}, h_{\text {prim }}^{p, q}$ be their respective dimensions. Lemma 8.13 yields

$$
\begin{align*}
\mathscr{H}^{p, q}(X, \mathbb{C}) & =\bigoplus_{r \geqslant(p+q-n)_{+}} L^{r} \operatorname{Prim} \mathscr{H}^{p-r, q-r}(X, \mathbb{C}),  \tag{8.14}\\
h^{p, q} & =\sum_{r \geqslant(p+q-n)_{+}} h_{\text {prim }}^{p-r, q-r} . \tag{8.15}
\end{align*}
$$

Formula (8.15) can be rewritten

$$
\left\{\begin{array}{l}
\text { if } p+q \leqslant n, \quad h^{p, q}=h_{\text {prim }}^{p, q}+h_{\text {prim }}^{p-1, q-1}+\cdots \\
\text { if } p+q \geqslant n, \quad h^{p, q}=h_{\text {prim }}^{n-q, n-p}+h_{\text {prim }}^{n-q-1, n-p-1}+\cdots .
\end{array}\right.
$$

(8.16) Corollary. The Hodge and Betti numbers satisfy the inequalities
a) if $k=p+q \leqslant n$, then $h^{p, q} \geqslant h^{p-1, q-1}, \quad b_{k} \geqslant b_{k-2}$,
b) if $k=p+q \geqslant n$, then $h^{p, q} \geqslant h^{p+1, q+1}, \quad b_{k} \geqslant b_{k+2}$.

Another important result of Hodge theory (which is in fact a direct consequence of Cor. 5.17) is the
(8.17) Hard Lefschetz theorem. The mappings

$$
\begin{array}{lcl}
L^{n-k}: & H^{k}(X, \mathbb{C}) \longrightarrow H^{2 n-k}(X, \mathbb{C}), & k \leqslant n, \\
L^{n-p-q}: & H^{p, q}(X, \mathbb{C}) \longrightarrow H^{n-q, n-p}(X, \mathbb{C}), & p+q \leqslant n,
\end{array}
$$

are isomorphisms.

## § 9. Jacobian and Albanese Varieties

## § 9.1. Description of the Picard Group

An important application of Hodge theory is a description of the Picard group $H^{1}\left(X, \mathscr{O}^{\star}\right)$ of a compact Kähler manifold. We assume here that $X$ is connected. The exponential exact sequence $0 \rightarrow \mathbb{Z} \rightarrow \mathbb{O} \rightarrow \mathbb{O}^{\star} \rightarrow 1$ gives

$$
\begin{align*}
0  \tag{9.1}\\
\quad \longrightarrow
\end{align*} H^{1}(X, \mathbb{Z}) \longrightarrow H^{1}(X, \mathbb{C}) \longrightarrow H^{1}\left(X, \mathbb{Q}^{\star}\right) ~ 子 H^{2}(X, \mathbb{Z}) \longrightarrow H^{2}(X, \mathbb{O}) \text { c. }
$$

because the map $\exp (2 \pi i \bullet): H^{0}(X, \mathscr{O})=\mathbb{C} \longrightarrow H^{0}\left(X, \mathscr{O}^{\star}\right)=\mathbb{C}^{\star}$ is onto. We have $H^{1}(X, \mathscr{\odot}) \simeq H^{0,1}(X, \mathbb{C})$ by $(\mathrm{V}-11.6)$. The dimension of this group is called the irregularity of $X$ and is usually denoted

$$
\begin{equation*}
q=q(X)=h^{0,1}=h^{1,0} . \tag{9.2}
\end{equation*}
$$

Therefore we have $b_{1}=2 q$ and

$$
\begin{equation*}
H^{1}(X, \mathbb{O}) \simeq \mathbb{C}^{q}, \quad H^{0}\left(X, \Omega_{X}^{1}\right)=H^{1,0}(X, \mathbb{C}) \simeq \mathbb{C}^{q} \tag{9.3}
\end{equation*}
$$

(9.4) Lemma. The image of $H^{1}(X, \mathbb{Z})$ in $H^{1}(X, \mathbb{O})$ is a lattice.

Proof. Consider the morphisms

$$
H^{1}(X, \mathbb{Z}) \longrightarrow H^{1}(X, \mathbb{R}) \longrightarrow H^{1}(X, \mathbb{C}) \longrightarrow H^{1}(X, \mathscr{O})
$$

induced by the inclusions $\mathbb{Z} \subset \mathbb{R} \subset \mathbb{C} \subset \mathbb{O}$. Since the Čech cohomology groups with values in $\mathbb{Z}, \mathbb{R}$ can be computed by finite acyclic coverings, we see that $H^{1}(X, \mathbb{Z})$ is a finitely generated $\mathbb{Z}$-module and that the image of $H^{1}(X, \mathbb{Z})$ in $H^{1}(X, \mathbb{R})$ is a lattice. It is enough to check that the map $H^{1}(X, \mathbb{R}) \longrightarrow H^{1}(X, \mathbb{O})$ is an isomorphism. However, the commutative diagram

shows that the map $H^{1}(X, \mathbb{R}) \longrightarrow H^{1}(X, \mathbb{O})$ corresponds in De Rham and Dolbeault cohomologies to the composite mapping

$$
H_{D R}^{1}(X, \mathbb{R}) \subset H_{D R}^{1}(X, \mathbb{C}) \longrightarrow H^{0,1}(X, \mathbb{C})
$$

Since $H^{1,0}(X, \mathbb{C})$ and $H^{0,1}(X, \mathbb{C})$ are complex conjugate subspaces in $H_{D R}^{1}(X, \mathbb{C})$, we conclude that $H_{D R}^{1}(X, \mathbb{R}) \longrightarrow H^{0,1}(X, \mathbb{C})$ is an isomorphism.

As a consequence of this lemma, $H^{1}(X, \mathbb{Z}) \simeq \mathbb{Z}^{2 q}$. The $q$-dimensional complex torus

$$
\begin{equation*}
\operatorname{Jac}(X)=H^{1}(X, \overparen{O}) / H^{1}(X, \mathbb{Z}) \tag{9.5}
\end{equation*}
$$

is called the Jacobian variety of $X$ and is isomorphic to the subgroup of $H^{1}\left(X, \mathbb{O}^{\star}\right)$ corresponding to line bundles of zero first Chern class. On the other hand, the kernel of

$$
H^{2}(X, \mathbb{Z}) \longrightarrow H^{2}(X, \mathbb{C})=H^{0,2}(X, \mathbb{C})
$$

which consists of integral cohomology classes of type $(1,1)$, is equal to the image of $c_{1}$ in $H^{2}(X, \mathbb{Z})$. This subgroup is called the Neron-Severi group of $X$, and is denoted $N S(X)$. The exact sequence (9.1) yields

$$
\begin{equation*}
0 \longrightarrow \mathrm{Jac}(X) \longrightarrow H^{1}\left(X, \overparen{O}^{\star}\right) \xrightarrow{c_{1}} N S(X) \longrightarrow 0 \tag{9.6}
\end{equation*}
$$

The Picard group $H^{1}\left(X, \mathscr{O}^{\star}\right)$ is thus an extension of the complex torus $\operatorname{Jac}(X)$ by the finitely generated $\mathbb{Z}$-module $N S(X)$.
(9.7) Corollary. The Picard group of $\mathbb{P}^{n}$ is $H^{1}\left(\mathbb{P}^{n}, \widehat{O}^{\star}\right) \simeq \mathbb{Z}$, and every line bundle over $\mathbb{P}^{n}$ is isomorphic to one of the line bundles $\mathbb{O}(k), k \in \mathbb{Z}$.

Proof. We have $H^{k}\left(\mathbb{P}^{n}, \mathbb{C}\right)=H^{0, k}\left(\mathbb{P}^{n}, \mathbb{C}\right)=0$ for $k \geqslant 1$ by Appl. 8.8, thus $\operatorname{Jac}\left(\mathbb{P}^{n}\right)=0$ and $N S\left(\mathbb{P}^{n}\right)=H^{2}\left(\mathbb{P}^{n}, \mathbb{Z}\right) \simeq \mathbb{Z}$. Moreover, $c_{1}(\mathbb{O}(1))$ is a generator of $H^{2}\left(\mathbb{P}^{n}, \mathbb{Z}\right)$ in virtue of Th. V-15.10.

## § 9.2. Albanese Variety

A proof similar to that of Lemma 9.4 shows that the image of $H^{2 n-1}(X, \mathbb{Z})$ in $H^{n-1, n}(X, \mathbb{C})$ via the composite map

$$
\begin{equation*}
H^{2 n-1}(X, \mathbb{Z}) \rightarrow H^{2 n-1}(X, \mathbb{R}) \rightarrow H^{2 n-1}(X, \mathbb{C}) \rightarrow H^{n-1, n}(X, \mathbb{C}) \tag{9.8}
\end{equation*}
$$

is a lattice. The $q$-dimensional complex torus

$$
\begin{equation*}
\operatorname{Alb}(X)=H^{n-1, n}(X, \mathbb{C}) / \operatorname{Im} H^{2 n-1}(X, \mathbb{Z}) \tag{9.9}
\end{equation*}
$$

is called the Albanese variety of $X$. We first give a slightly different description of $\operatorname{Alb}(X)$, based on the Serre duality isomorphism

$$
H^{n-1, n}(X, \mathbb{C}) \simeq\left(H^{1,0}(X, \mathbb{C})\right)^{\star} \simeq\left(H^{0}\left(X, \Omega_{X}^{1}\right)\right)^{\star}
$$

(9.10) Lemma. For any compact oriented differentiable manifold $M$ with $\operatorname{dim}_{\mathbb{R}} M=m$, there is a natural isomorphism

$$
H_{1}(M, \mathbb{Z}) \rightarrow H^{m-1}(M, \mathbb{Z})
$$

where $H_{1}(M, \mathbb{Z})$ is the first homology group of $M$, that is, the abelianization of $\pi_{1}(M)$.
Proof. This is a well known consequence of Poincaré duality, see e.g. [Spanier 1966]. We will content ourselves with a description of the morphism. Fix a base point $a \in M$. Every homotopy class $[\gamma] \in \pi_{1}(M, a)$ can be represented by as a composition of closed loops diffeomorphic to $S^{1}$. Let $\gamma$ be such a loop. As every oriented vector bundle over $S^{1}$ is trivial, the normal bundle to $\gamma$ is trivial. Hence $\gamma\left(S^{1}\right)$ has a neighborhood $U$ diffeomorphic to $S^{1} \times \mathbb{R}^{m-1}$, and there is a diffeomorphism $\varphi: S^{1} \times \mathbb{R}^{m-1} \rightarrow U$ with $\varphi_{\mid S^{1} \times\{0\}}=\gamma$. Let $\left\{\delta_{0}\right\} \in H_{c}^{m-1}\left(\mathbb{R}^{m-1}, \mathbb{Z}\right)$ be the fundamental class represented by the Dirac measure $\delta_{0} \in \mathscr{D}_{0}^{\prime}\left(\mathbb{R}^{m-1}\right)$ in De Rham cohomology. Then the cartesian product $1 \times$ $\left\{\delta_{0}\right\} \in H_{c}^{m-1}\left(S^{1} \times \mathbb{R}^{m-1}, \mathbb{Z}\right)$ is represented by the current $\left[S^{1}\right] \otimes\left\{\delta_{0}\right\} \in \mathscr{D}_{1}^{\prime}\left(S^{1} \times \mathbb{R}^{m-1}\right)$ and the current of integration over $\gamma$ is precisely the direct image current

$$
I_{\gamma}:=\varphi_{\star}\left(\left[S^{1}\right] \otimes \delta_{0}\right)=\left(\varphi^{-1}\right)^{\star}\left(\left[S^{1}\right] \otimes \delta_{0}\right) .
$$

Its cohomology class $\left\{I_{\gamma}\right\} \in H_{c}^{m-1}(U, \mathbb{R})$ is thus the image of the class

$$
\left(\varphi^{-1}\right)^{\star}\left(1 \times\left\{\delta_{0}\right\}\right) \in H_{c}^{m-1}(U, \mathbb{Z}) .
$$

Therefore, we have obtained a well defined morphism

$$
\pi_{1}(M, a) \longrightarrow H_{c}^{m-1}(U, \mathbb{Z}) \longrightarrow H^{m-1}(M, \mathbb{Z}), \quad[\gamma] \longmapsto\left(\varphi^{-1}\right)^{\star}\left(1 \times\left\{\delta_{0}\right\}\right)
$$

and the image of $[\gamma]$ in $H^{m-1}(M, \mathbb{R})$ is the De Rham cohomology class of the integration current $I_{\gamma}$.

Thanks to Lemma 9.10, we can reformulate the definition of the Albanese variety as

$$
\begin{equation*}
\operatorname{Alb}(X)=\left(H^{0}\left(X, \Omega_{X}^{1}\right)\right)^{\star} / \operatorname{Im} H_{1}(X, \mathbb{Z}) \tag{9.11}
\end{equation*}
$$

where $H_{1}(X, \mathbb{Z})$ is mapped to $\left(H^{0}\left(X, \Omega_{X}^{1}\right)\right)^{\star}$ by

$$
[\gamma] \longmapsto \widetilde{I}_{\gamma}=\left(u \mapsto \int_{\gamma} u\right) .
$$

Observe that the integral only depends on the homotopy class $[\gamma]$ because all holomorphic 1 -forms $u$ on $X$ are closed by Prop. 8.9.

We are going to show that there exists a canonical holomorphic map $\alpha: X \rightarrow \operatorname{Alb}(X)$. Let $a$ be a base point in $X$. For any $x \in X$, we select a path $\xi$ from $a$ to $x$ and associate to $x$ the linear form in $\left(H^{0}\left(X, \Omega_{X}^{1}\right)\right)^{\star}$ defined by $\widetilde{I}_{\xi}$. By construction the class of this linear form $\bmod \operatorname{Im} H_{1}(X, \mathbb{Z})$ does not depend on $\xi$, since $\widetilde{I}_{\xi^{\prime}-1}$ is in the image of $H_{1}(X, \mathbb{Z})$ for any other path $\xi^{\prime}$. It is thus legitimate to define the Albanese map as

$$
\begin{equation*}
\alpha: X \longrightarrow \operatorname{Alb}(X), \quad x \longmapsto\left(u \mapsto \int_{a}^{x} u\right) \bmod \operatorname{Im} H_{1}(X, \mathbb{Z}) \tag{9.12}
\end{equation*}
$$

Of course, if we fix a basis $\left(u_{1}, \ldots, u_{q}\right)$ of $H^{0}\left(X, \Omega_{X}^{1}\right)$, the Albanese map can be seen in coordinates as the map

$$
\begin{equation*}
\alpha: X \longrightarrow \mathbb{C}^{q} / \Lambda, \quad x \longmapsto\left(\int_{a}^{x} u_{1}, \ldots, \int_{a}^{x} u_{q}\right) \bmod \Lambda, \tag{9.13}
\end{equation*}
$$

where $\Lambda \subset \mathbb{C}^{q}$ is the group of periods of $\left(u_{1}, \ldots, u_{q}\right)$ :

$$
\Lambda=\left\{\left(\int_{\gamma} u_{1}, \ldots, \int_{\gamma} u_{q}\right) ;[\gamma] \in \pi_{1}(X, a)\right\} .
$$

It is then clear that $\alpha$ is a holomorphic map. With the original definition (9.9) of the Albanese variety, it is not difficult to see that $\alpha$ is the map given by

$$
\begin{equation*}
\alpha: X \longrightarrow \operatorname{Alb}(X), \quad x \longmapsto\left\{I_{\xi}^{n-1, n}\right\} \bmod H^{2 n-1}(X, \mathbb{Z}) \tag{9.14}
\end{equation*}
$$

where $\left\{I_{\xi}^{n-1, n}\right\} \in H^{n-1, n}(X, \mathbb{C})$ denotes the $(n-1, n)$-component of the De Rham cohomology class $\left\{I_{\xi}\right\}$.

## §10. Complex Curves

We show here how Hodge theory can be used to derive quickly the basic properties of compact manifolds of complex dimension 1 (also called complex curves or Riemann surfaces). Let $X$ be such a curve. We shall always assume in this section that $X$ is compact and connected. Since every positive (1,1)-form on a curve defines a Kähler metric, the results of $\S 8$ and $\S 9$ can be applied.

## $\S$ 10.1. Riemann-Roch Formula

Denoting $g=h^{1}(X, \overparen{O})$, we find

$$
\begin{align*}
& H^{1}(X, \mathbb{O}) \simeq \mathbb{C}^{g}, \quad H^{0}\left(X, \Omega_{X}^{1}\right) \simeq \mathbb{C}^{g}  \tag{10.1}\\
& H^{0}(X, \mathbb{Z})=\mathbb{Z}, \quad H^{1}(X, \mathbb{Z})=\mathbb{Z}^{2 g}, \quad H^{2}(X, \mathbb{Z})=\mathbb{Z} \tag{10.2}
\end{align*}
$$

The classification of oriented topological surfaces shows that $X$ is homeomorphic to a sphere with $g$ handles ( $=$ torus with $g$ holes), but this property will not be used in the sequel. The number $g$ is called the genus of $X$.

Any divisor on $X$ can be written $\Delta=\sum m_{j} a_{j}$ where $\left(a_{j}\right)$ is a finite sequence of points and $m_{j} \in \mathbb{Z}$. Let $E$ be a line bundle over $X$. We shall identify $E$ and the associated locally free sheaf $\mathscr{O}(E)$. According to V-13.2, we denote by $E(\Delta)$ the sheaf of germs of meromorphic sections $f$ of $E$ such that $\operatorname{div} f+\Delta \geqslant 0$, i.e. which have a pole of order $\leqslant m_{j}$ at $a_{j}$ if $m_{j}>0$, and which have a zero of order $\geqslant\left|m_{j}\right|$ at $a_{j}$ if $m_{j}<0$. Clearly

$$
\begin{equation*}
E(\Delta)=E \otimes \mathscr{O}(\Delta), \quad \mathscr{O}\left(\Delta+\Delta^{\prime}\right)=\mathscr{O}(\Delta) \otimes \mathscr{O}\left(\Delta^{\prime}\right) \tag{10.3}
\end{equation*}
$$

For any point $a \in X$ and any integer $m>0$, there is an exact sequence

$$
0 \longrightarrow E \longrightarrow E(m[a]) \longrightarrow \mathscr{S} \longrightarrow 0
$$

where $\mathscr{S}=E(m[a]) / E$ is a sheaf with only one non zero stalk $\mathscr{S}_{a}$ isomorphic to $\mathbb{C}^{m}$. Indeed, if $z$ is a holomorphic coordinate near $a$, the stalk $\mathscr{S}_{a}$ corresponds to the polar parts $\sum_{-m \leqslant k<0} c_{k} z^{k}$ in the power series expansions of germs of meromorphic sections at point $a$. We get an exact sequence

$$
H^{0}(X, E(m[a])) \longrightarrow \mathbb{C}^{m} \longrightarrow H^{1}(X, E)
$$

When $m$ is chosen larger than $\operatorname{dim} H^{1}(X, E)$, we see that $E(m[a])$ has a non zero section and conclude:
(10.4) Theorem. Let a be a given point on a curve. Then every line bundle $E$ has non zero meromorphic sections $f$ with a pole at a and no other poles.

If $\Delta$ is the divisor of a meromorphic section $f$ of $E$, we have $E \simeq \mathscr{O}(\Delta)$, so the map

$$
\operatorname{Div}(X) \longrightarrow H^{1}\left(X, \mathscr{O}^{\star}\right), \quad \Delta \longmapsto \mathscr{O}(\Delta)
$$

is onto (cf. (V-13.8)). On the other hand, Div is clearly a soft sheaf, thus $H^{1}(X, \operatorname{Div})=0$. The long cohomology sequence associated to the exact sequence $1 \rightarrow \mathscr{O}^{\star} \rightarrow \mathcal{M}^{\star} \rightarrow \mathrm{Div} \rightarrow$ 0 implies:
(10.5) Corollary. On any complex curve, one has $H^{1}\left(X, \mathcal{M}^{\star}\right)=0$ and there is an exact sequence

$$
0 \longrightarrow \mathbb{C}^{\star} \longrightarrow \mathcal{M}^{\star}(X) \longrightarrow \operatorname{Div}(X) \longrightarrow H^{1}\left(X, \mathscr{O}^{\star}\right) \longrightarrow 0
$$

The first Chern class $c_{1}(E) \in H^{2}(X, \mathbb{Z})$ can be interpreted as an integer. This integer is called the degree of $E$. If $E \simeq \mathscr{O}(\Delta)$ with $\Delta=\sum m_{j} a_{j}$, formula V-13.6 shows that the image of $c_{1}(E)$ in $H^{2}(X, \mathbb{R})$ is the De Rham cohomology class of the associated current $[\Delta]=\sum m_{j} \delta_{a_{j}}$, hence

$$
\begin{equation*}
c_{1}(E)=\int_{X}[\Delta]=\sum m_{j} . \tag{10.6}
\end{equation*}
$$

If $\sum m_{j} a_{j}$ is the divisor of a meromorphic function, we have $\sum m_{j}=0$ because the associated bundle $E=\mathscr{O}\left(\sum m_{j} a_{j}\right)$ is trivial.
(10.7) Theorem. Let $E$ be a line bundle on a complex curve $X$. Then
a) $H^{0}(X, E)=0$ if $c_{1}(E)<0$ or if $c_{1}(E)=0$ and $E$ is non trivial;
b) For every positive $(1,1)$-form $\omega$ on $X$ with $\int_{X} \omega=1$, $E$ has a hermitian metric such that $\frac{\mathrm{i}}{2 \pi} \Theta(E)=c_{1}(E) \omega$. In particular, $E$ has a metric of positive (resp. negative) curvature if and only if $c_{1}(E)>0$ (resp. if and only if $\left.c_{1}(E)<0\right)$.

Proof. a) If $E$ has a non zero holomorphic section $f$, then its degree is $c_{1}(E)=\int_{X} \operatorname{div} f \geqslant$ 0 . In fact, we even have $c_{1}(E)>0$ unless $f$ does not vanish, in which case $E$ is trivial.
b) Select an arbitrary hermitian metric $h$ on $E$. Then $c_{1}(E) \omega-\frac{i}{2 \pi} \Theta_{h}(E)$ is a real $(1,1)$-form cohomologous to zero (the integral over $X$ is zero), so Lemma 8.6 c) implies

$$
c_{1}(E) \omega-\frac{\mathrm{i}}{2 \pi} \Theta_{h}(E)=\mathrm{i} d^{\prime} d^{\prime \prime} \varphi
$$

for some real function $\varphi \in \mathscr{C}^{\infty}(X, \mathbb{R})$. If we replace the initial metric of $E$ by $h^{\prime}=h e^{-\varphi}$, we get a metric of constant curvature $c_{1}(E) \omega$.
(10.8) Riemann-Roch formula. Let $E$ be a holomorphic line bundle and let $h^{q}(E)=$ $\operatorname{dim} H^{q}(X, E)$. Then

$$
h^{0}(E)-h^{1}(E)=c_{1}(E)-g+1
$$

Moreover $h^{1}(E)=h^{0}\left(K \otimes E^{\star}\right)$, where $K=\Omega_{X}^{1}$ is the canonical line bundle of $X$.
Proof. We claim that for every line bundle $F$ and every divisor $\Delta$ we have the equality

$$
\begin{equation*}
h^{0}(F(\Delta))-h^{1}(F(\Delta))=h^{0}(F)-h^{1}(F)+\int_{X}[\Delta] . \tag{10.9}
\end{equation*}
$$

If we write $E=\mathscr{O}(\Delta)$ and apply the above equality with $F=\mathscr{O}$, the Riemann-Roch formula results from (10.6), (10.9) and from the equalities

$$
h^{0}(\mathbb{C})=\operatorname{dim} H^{0}(X, \mathbb{C})=1, \quad h^{1}(\mathbb{C})=\operatorname{dim} H^{1}(X, \mathbb{O})=g .
$$

However, (10.9) need only be proved when $\Delta \geqslant 0$ : otherwise we are reduced to this case by writing $\Delta=\Delta_{1}-\Delta_{2}$ with $\Delta_{1}, \Delta_{2} \geqslant 0$ and by applying the result to the pairs ( $F, \Delta_{1}$ ) and $\left(F(\Delta), \Delta_{2}\right)$. If $\Delta=\sum m_{j} a_{j} \geqslant 0$, there is an exact sequence

$$
0 \longrightarrow F \longrightarrow F(\Delta) \longrightarrow \mathscr{S} \longrightarrow 0
$$

where $\mathscr{S}_{a_{j}} \simeq \mathbb{C}^{m_{j}}$ and the other stalks are zero. Let $m=\sum m_{j}=\int_{X}[\Delta]$. The sheaf $\mathscr{S}$ is acyclic, because its support $\left\{a_{j}\right\}$ is of dimension 0 . Hence there is an exact sequence

$$
0 \longrightarrow H^{0}(F) \longrightarrow H^{0}(F(\Delta)) \longrightarrow \mathbb{C}^{m} \longrightarrow H^{1}(F) \longrightarrow H^{1}(F(\Delta)) \longrightarrow 0
$$

and (10.9) follows. The equality $h^{1}(E)=h^{0}\left(K \otimes E^{\star}\right)$ is a consequence of the Serre duality theorem

$$
\left(H^{0,1}(X, E)\right)^{\star} \simeq H^{1,0}\left(X, E^{\star}\right), \quad \text { i.e. }\left(H^{1}(X, E)\right)^{\star} \simeq H^{0}\left(X, K \otimes E^{\star}\right)
$$

(10.10) Corollary (Hurwitz' formula). $c_{1}(K)=2 g-2$.

Proof. Apply Riemann-Roch to $E=K$ and observe that

$$
\begin{align*}
& h^{0}(K)=\operatorname{dim} H^{0}\left(X, \Omega_{X}^{1}\right)=g \\
& h^{1}(K)=\operatorname{dim} H^{1}\left(X, \Omega_{X}^{1}\right)=h^{1,1}=b_{2}=1 \tag{10.11}
\end{align*}
$$

(10.12) Corollary. For every $a \in X$ and every $m \in \mathbb{Z}$

$$
h^{0}(K(-m[a]))=h^{1}(\mathbb{O}(m[a]))=h^{0}(\mathbb{O}(m[a]))-m+g-1 .
$$

## § 10.2. Jacobian of a Curve

By the Neron-Severi sequence (9.6), there is an exact sequence

$$
\begin{equation*}
0 \longrightarrow \mathrm{Jac}(X) \longrightarrow H^{1}\left(X, \mathbb{O}^{\star}\right) \xrightarrow{c_{1}} \mathbb{Z} \longrightarrow 0, \tag{10.13}
\end{equation*}
$$

where the $\operatorname{Jacobian} \operatorname{Jac}(X)$ is a $g$-dimensional torus. Choose a base point $a \in X$. For every point $x \in X$, the line bundle $\mathbb{O}([x]-[a])$ has zero first Chern class, so we have a well-defined map

$$
\begin{equation*}
\Phi_{a}: X \longrightarrow \operatorname{Jac}(X), \quad x \longmapsto \mathscr{O}([x]-[a]) . \tag{10.14}
\end{equation*}
$$

Observe that the $\operatorname{Jacobian} \operatorname{Jac}(X)$ of a curve coincides by definition with the Albanese variety $\operatorname{Alb}(X)$.
(10.15) Lemma. The above map $\Phi_{a}$ coincides with the Albanese map $\alpha: X \longrightarrow \operatorname{Alb}(X)$ defined in (9.12).

Proof. By holomorphic continuation, it is enough to prove that $\Phi_{a}(x)=\alpha(x)$ when $x$ is near $a$. Let $z$ be a complex coordinate and let $D^{\prime} \subset \subset D$ be open disks centered at $a$. Relatively to the covering

$$
U_{1}=D, \quad U_{2}=X \backslash \overline{D^{\prime}}
$$

the line bundle $\mathscr{O}([x]-[a])$ is defined by the Čech cocycle $c \in C^{1}\left(\mathscr{U}, \mathscr{O}^{\star}\right)$ such that

$$
c_{12}(z)=\frac{z-x}{z-a} \quad \text { on } \quad U_{12}=D \backslash \overline{D^{\prime}} .
$$

On the other hand, we compute $\alpha(x)$ by Formula (9.14). The path integral current $I_{[a, x]} \in$ $\mathscr{D}_{1}^{\prime}(X)$ is equal to 0 on $U_{2}$. Lemma I-2.10 implies $d^{\prime \prime}(d z / 2 \pi \mathrm{i} z)=\delta_{0} d \bar{z} \wedge d z / 2 i=\delta_{0}$ according to the usual identification of distributions and currents of degree 0 , thus

$$
I_{[a, x]}^{0,1}=d^{\prime \prime}\left(\frac{d z}{2 \pi i z} \star I_{[a, x]}^{0,1}\right) \quad \text { on } \quad U_{1} .
$$

Therefore $\left\{I_{[a, x]}^{0,1}\right\} \in H^{0,1}(X, \mathbb{C})$ is equal to the Čech cohomology class $\left\lceil^{\prime}\right\}$ in $H^{1}(X, \mathbb{O})$ represented by the cocycle

$$
c_{12}^{\prime}(z)=\frac{d w}{2 \pi i w} \star I_{[a, x]}^{0,1}=\frac{1}{2 \pi i} \int_{a}^{x} \frac{d w}{w-z}=\frac{1}{2 \pi \mathrm{i}} \log \frac{z-x}{z-a} \text { on } U_{12}
$$

and we have $c=\exp \left(2 \pi \mathrm{i} c^{\prime}\right)$ in $H^{1}\left(X, \mathscr{O}^{\star}\right)$.
The nature of $\Phi_{a}$ depends on the value of the genus $g$. A careful examination of $\Phi_{a}$ will enable us to determine all curves of genus 0 and 1 .
(10.16) Theorem. The following properties are equivalent:
a) $g=0$;
b) $X$ has a meromorphic function $f$ having only one simple pole $p$;
c) $X$ is biholomorphic to $\mathbb{P}^{1}$.

Proof. c) $\Longrightarrow$ a) is clear.
a) $\Longrightarrow \mathrm{b})$. Since $g=0$, we have $\operatorname{Jac}(X)=0$. If $p, p^{\prime} \in X$ are distinct points, the bundle $\mathbb{O}\left(\left[p^{\prime}\right]-[p]\right)$ has zero first Chern class, therefore it is trivial and there exists a meromorphic function $f$ with div $f=\left[p^{\prime}\right]-[p]$. In particular $p$ is the only pole of $f$, and this pole is simple.
b) $\Longrightarrow$ c). We may consider $f$ as a map $X \longrightarrow \mathbb{P}^{1}=\mathbb{C} \cup\{\infty\}$. For every value $w \in \mathbb{C}$, the function $f-w$ must have exactly one simple zero $x \in X$ because $\int_{X} \operatorname{div}(f-w)=0$ and $p$ is a simple pole. Therefore $f: X \rightarrow \mathbb{P}^{1}$ is bijective and $X$ is biholomorphic to $\mathbb{P}^{1}$.
(10.17) Theorem. The map $\Phi_{a}$ is always injective for $g \geqslant 1$.
a) If $g=1, \Phi_{a}$ is a biholomorphism. In particular every curve of genus 1 is biholomorphic to a complex torus $\mathbb{C} / \Gamma$.
b) If $g \geqslant 2, \Phi_{a}$ is an embedding.

Proof. If $\Phi_{a}$ is not injective, there exist points $x_{1} \neq x_{2}$ such that $\mathscr{O}\left(\left[x_{1}\right]-\left[x_{2}\right]\right)$ is trivial; then there is a meromorphic function $f$ such that div $f=\left[x_{1}\right]-\left[x_{2}\right]$ and Th. 10.16 implies that $g=0$.

When $g=1, \Phi_{a}$ is an injective map $X \longrightarrow \operatorname{Jac}(X) \simeq \mathbb{C} / \Gamma$, thus $\Phi_{a}$ is open. It follows that $\Phi_{a}(X)$ is a compact open subset of $\mathbb{C} / \Gamma$, so $\Phi_{a}(X)=\mathbb{C} / \Gamma$ and $\Phi_{a}$ is a biholomorphism of $X$ onto $\mathbb{C} / \Gamma$.

In order to prove that $\Phi_{a}$ is an embedding when $g \geqslant 2$, it is sufficient to show that the holomorphic 1-forms $u_{1}, \ldots, u_{g}$ do not all vanish at a given point $x \in X$. In fact, $X$ has no non constant meromorphic function having only a simple pole at $x$, thus $h^{0}(\mathbb{O}([x]))=1$ and Cor. 10.12 implies

$$
h^{0}(K(-[x]))=g-1<h^{0}(K)=g .
$$

Hence $K$ has a section $u$ which does not vanish at $x$.

## $\S$ 10.3. Weierstrass Points of a Curve

We want to study how many meromorphic functions have a unique pole of multiplicity $\leqslant m$ at a given point $a \in X$, i.e. we want to compute $h^{0}(\mathscr{O}(m[a]))$. As we shall see soon, these numbers may depend on $a$ only if $m$ is small. We have $c_{1}(K(-m[a]))=2 g-2-m$, so the degree is $<0$ and $h^{0}(K(-m[a]))=0$ for $m \geqslant 2 g-1$ by 10.7 a). Cor. 10.12 implies

$$
\begin{equation*}
h^{0}(\mathbb{O}(m[a]))=m-g+1 \quad \text { for } \quad m \geqslant 2 g-1 . \tag{10.18}
\end{equation*}
$$

It remains to compute $h^{0}(K(-m[a]))$ for $0 \leqslant m \leqslant 2 g-2$ and $g \geqslant 1$. Let $u_{1}, \ldots, u_{g}$ be a basis of $H^{0}(X, K)$ and let $z$ be a complex coordinate centered at $a$. Any germ $u \in \mathbb{O}(K)_{a}$ can be written $u=U(z) d z$ with $U(z)=\sum_{m \in \mathbb{N}} \frac{1}{m!} U^{(m)}(a) z^{m} d z$. The unique non zero stalk of the quotient sheaf $\mathbb{O}(K(-m[a])) / \mathscr{O}(K(-(m+1)[a]))$ is canonically isomorphic to $K_{a}^{m+1}$ via the map $u \mapsto U^{(m)}(a)(d z)^{m+1}$, which is independant of the choice of $z$. Hence $\bigwedge^{g}(\mathbb{O}(K) / \mathscr{O}(K-g[a])) \simeq K_{a}^{1+2+\ldots+g}$ and the Wronskian

$$
W\left(u_{1}, \ldots, u_{g}\right)=\left|\begin{array}{ccc}
U_{1}(z) & \ldots & U_{g}(z)  \tag{10.19}\\
U_{1}^{\prime}(z) & \ldots & U_{g}^{\prime}(z) \\
\vdots & & \vdots \\
U_{1}^{(g-1)}(z) & \ldots & U_{g}^{(g-1)}(z)
\end{array}\right| d z^{1+2+\ldots+g}
$$

defines a global section $W\left(u_{1}, \ldots, u_{g}\right) \in H^{0}\left(X, K^{g(g+1) / 2}\right)$. At the given point $a$, we can find linear combinations $\widetilde{u}_{1}, \ldots, \widetilde{u}_{g}$ of $u_{1}, \ldots, u_{g}$ such that

$$
\widetilde{u}_{j}(z)=\left(z^{s_{j}-1}+\mathrm{O}\left(z^{s_{j}}\right)\right) d z, \quad s_{1}<\ldots<s_{g} .
$$

We know that not all sections of $K$ vanish at $a$ and that $c_{1}(K)=2 g-2$, thus $s_{1}=1$ and $s_{g} \leqslant 2 g-1$. We have $W\left(\widetilde{u}_{1}, \ldots, \widetilde{u}_{g}\right) \sim W\left(z^{s_{1}-1} d z, \ldots, z^{s_{g}-1} d z\right)$ at point $a$, and an easy induction on $\sum s_{j}$ combined with differentiation in $z$ yields

$$
W\left(z^{s_{1}-1} d z, \ldots, z^{s_{g}-1} d z\right)=C z^{s_{1}+\ldots+s_{g}-g(g+1) / 2} d z^{g(g+1) / 2}
$$

for some positive integer constant $C$. In particular, $W\left(u_{1}, \ldots, u_{g}\right)$ is not identically zero and vanishes at $a$ with multiplicity

$$
\begin{equation*}
\mu_{a}=s_{1}+\ldots+s_{g}-g(g+1) / 2>0 \tag{10.20}
\end{equation*}
$$

unless $s_{1}=1, s_{2}=2, \ldots, s_{g}=g$. Now, we have

$$
h^{0}(K(-m[a]))=\operatorname{card}\left\{j ; s_{j}>m\right\}=g-\operatorname{card}\left\{j ; s_{j} \leqslant m\right\}
$$

and Cor. 10.12 gives

$$
\begin{equation*}
h^{0}(\mathbb{O}(m[a]))=m+1-\operatorname{card}\left\{j ; s_{j} \leqslant m\right\} . \tag{10.21}
\end{equation*}
$$

If $a$ is not a zero of $W\left(u_{1}, \ldots, u_{g}\right)$, we find

$$
\begin{cases}h^{0}(\mathbb{O}(m[a]))=1 & \text { for } m \leqslant g,  \tag{10.22}\\ h^{0}(\mathbb{O}(m[a]))=m+1-g & \text { for } m>g .\end{cases}
$$

The zeroes of $W\left(u_{1}, \ldots, u_{g}\right)$ are called the Weierstrass points of $X$, and the associated Weierstrass sequence is the sequence $w_{m}=h^{0}(\mathbb{O}(m[a])), m \in \mathbb{N}$. We have $w_{m-1} \leqslant w_{m} \leqslant$ $w_{m-1}+1$ and $s_{1}<\ldots<s_{g}$ are precisely the integers $m \geqslant 1$ such that $w_{m}=w_{m-1}$. The numbers $s_{j} \in\{1,2, \ldots, 2 g-1\}$ are called the gaps and $\mu_{a}$ the weight of the Weierstrass point $a$. Since $W\left(u_{1}, \ldots, u_{g}\right)$ is a section of $K^{g(g+1) / 2}$, Hurwitz' formula implies

$$
\begin{equation*}
\sum_{a \in X} \mu_{a}=c_{1}\left(K^{g(g+1) / 2}\right)=g(g+1)(g-1) . \tag{10.23}
\end{equation*}
$$

In particular, a curve of genus $g$ has at most $g(g+1)(g-1)$ Weierstrass points.

## § 11. Hodge-Frölicher Spectral Sequence

Let $X$ be a compact complex $n$-dimensional manifold. We consider the double complex $K^{p, q}=\mathscr{C}^{\infty}\left(X, \Lambda^{p, q} T_{X}^{\star}\right), d=d^{\prime}+d^{\prime \prime}$. The Hodge-Frölicher spectral sequence is by definition the spectral sequence associated to this double complex (cf. IV-11.9). It starts with

$$
\begin{equation*}
E_{1}^{p, q}=H^{p, q}(X, \mathbb{C}) \tag{11.1}
\end{equation*}
$$

and the limit term $E_{\infty}^{p, q}$ is the graded module associated to a filtration of the De Rham cohomology group $H^{k}(X, \mathbb{C}), k=p+q$. In particular, if the numbers $b_{k}$ and $h^{p, q}$ are still defined as in (8.11), we have

$$
\begin{equation*}
b_{k}=\sum_{p+q=k} \operatorname{dim} E_{\infty}^{p, q} \leqslant \sum_{p+q=k} \operatorname{dim} E_{1}^{p, q}=\sum_{p+q=k} h^{p, q} . \tag{11.2}
\end{equation*}
$$

The equality is equivalent to the degeneration of the spectral sequence at $E_{1}^{\bullet}$. As a consequence, the Hodge-Frölicher spectral sequence of a compact Kähler manifold degenerates in $E_{1}^{\bullet}$.
(11.3) Theorem and Definition. The existence of an isomorphism

$$
H_{\mathrm{DR}}^{k}(X, \mathbb{C}) \simeq \bigoplus_{p+q=k} H^{p, q}(X, \mathbb{C})
$$

is equivalent to the degeneration of the Hodge-Frölicher spectral sequence at $E_{1}$. In this case, the isomorphism is canonically defined and we say that $X$ admits a Hodge decomposition.

In general, interesting informations can be deduced from the spectral sequence. Theorem IV-11.8 shows in particular that

$$
\begin{equation*}
b_{1} \geqslant \operatorname{dim} E_{2}^{1,0}+\left(\operatorname{dim} E_{2}^{0,1}-\operatorname{dim} E_{2}^{2,0}\right)_{+} . \tag{11.4}
\end{equation*}
$$

However, $E_{2}^{1,0}$ is the central cohomology group in the sequence

$$
d_{1}=d^{\prime}: E_{1}^{0,0} \longrightarrow E_{1}^{1,0} \longrightarrow E_{1}^{2,0}
$$

and as $E_{1}^{0,0}$ is the space of holomorphic functions on $X$, the first map $d_{1}$ is zero (by the maximum principle, holomorphic functions are constant on each connected component of $X)$. Hence $\operatorname{dim} E_{2}^{1,0} \geqslant h^{1,0}-h^{2,0}$. Similarly, $E_{2}^{0,1}$ is the kernel of a map $E_{1}^{0,1} \rightarrow E_{1}^{1,1}$, thus $\operatorname{dim} E_{2}^{0,1} \geqslant h^{0,1}-h^{1,1}$. By (11.4) we obtain

$$
\begin{equation*}
b_{1} \geqslant\left(h^{1,0}-h^{2,0}\right)_{+}+\left(h^{0,1}-h^{1,1}-h^{2,0}\right)_{+} . \tag{11.5}
\end{equation*}
$$

Another interesting relation concerns the topological Euler-Poincaré characteristic

$$
\chi_{\mathrm{top}}(X)=b_{0}-b_{1}+\ldots-b_{2 n-1}+b_{2 n}
$$

We need the following simple lemma.
(11.6) Lemma. Let $\left(C^{\bullet}, d\right)$ a bounded complex of finite dimensional vector spaces over some field. Then, the Euler characteristic

$$
\chi\left(C^{\bullet}\right)=\sum(-1)^{q} \operatorname{dim} C^{q}
$$

is equal to the Euler characteristic $\chi\left(H^{\bullet}\left(C^{\bullet}\right)\right)$ of the cohomology module.
Proof. Set

$$
c_{q}=\operatorname{dim} C^{q}, \quad z_{q}=\operatorname{dim} Z^{q}\left(C^{\bullet}\right), \quad b_{q}=\operatorname{dim} B^{q}\left(C^{\bullet}\right), \quad h_{q}=\operatorname{dim} H^{q}\left(C^{\bullet}\right)
$$

Then

$$
c_{q}=z_{q}+b_{q+1}, \quad h_{q}=z_{q}-b_{q} .
$$

Therefore we find

$$
\sum(-1)^{q} c_{q}=\sum(-1)^{q} z_{q}-\sum(-1)^{q} b_{q}=\sum(-1)^{q} h_{q}
$$

In particular, if the term $E_{r}^{\bullet}$ of the spectral sequence of a filtered complex $K^{\bullet}$ is a bounded and finite dimensional complex, we have

$$
\chi\left(E_{r}^{\bullet}\right)=\chi\left(E_{r+1}^{\bullet}\right)=\ldots=\chi\left(E_{\infty}^{\bullet}\right)=\chi\left(H^{\bullet}\left(K_{\bullet}^{\bullet}\right)\right)
$$

because $E_{r+1}^{\bullet}=H^{\bullet}\left(E_{r}^{\bullet}\right)$ and $\operatorname{dim} E_{\infty}^{l}=\operatorname{dim} H^{l}\left(K^{\bullet}\right)$. In the Hodge-Frölicher spectral sequence, we have $\operatorname{dim} E_{1}^{l}=\sum_{p+q=l} h^{p, q}$, hence:
(11.7) Theorem. For any compact complex manifold $X$, one has

$$
\chi_{\mathrm{top}}(X)=\sum_{0 \leqslant k \leqslant 2 n}(-1)^{k} b_{k}=\sum_{0 \leqslant p, q \leqslant n}(-1)^{p+q} h^{p, q} .
$$

## § 12. Effect of a Modification on Hodge Decomposition

In this section, we show that the existence of a Hodge decomposition on a compact complex manifold $X$ is guaranteed as soon as there exists such a decomposition on a modification $\widetilde{X}$ of $X$ (see II-??.?? for the Definition). This leads us to extend Hodge theory to a class of manifolds which are non necessarily Kähler, the so called Fujiki class ( $C)$ ) of manifolds bimeromorphic to Kähler manifolds.

## $\S$ 12.1. Sheaf Cohomology Reinterpretation of $H_{B C}^{p, q}(X, \mathbb{C})$

We first give a description of $H_{\mathrm{BC}}^{p, q}(X, \mathbb{C})$ in terms of the hypercohomology of a suitable complex of sheaves. This interpretation, combined with the analogue of the HodgeFrölicher spectral sequence, will imply in particular that $H_{\mathrm{BC}}^{p, q}(X, \mathbb{C})$ is always finite dimensional when $X$ is compact. Let us denote by $\mathscr{E}^{p, q}$ the sheaf of germs of $\mathscr{C}^{\infty}$ forms of bidegree $(p, q)$, and by $\Omega^{p}$ the sheaf of germs of holomorphic $p$-forms on $X$. For a fixed bidegree $\left(p_{0}, q_{0}\right)$, we let $k_{0}=p_{0}+q_{0}$ and we introduce a complex of sheaves ( $\left.\mathscr{L}_{p_{0}, q_{0}}^{\bullet}, \delta\right)$, also denoted $\mathscr{L}^{\bullet}$ for simplicity, such that

$$
\begin{aligned}
\mathscr{L}^{k} & =\bigoplus_{p+q=k, p<p_{0}, q<q_{0}} \mathscr{E}^{p, q} \quad \text { for } k \leqslant k_{0}-2, \\
\mathscr{L}^{k-1} & =\bigoplus_{p+q=k, p \geqslant p_{0}, q \geqslant q_{0}} \mathscr{E}^{p, q} \quad \text { for } k \geqslant k_{0} .
\end{aligned}
$$

The differential $\delta^{k}$ on $\mathscr{L}^{k}$ is chosen equal to the exterior derivative $d$ for $k \neq k_{0}-2$ (in the case $k \leqslant k_{0}-3$, we neglect the components which fall outside $\mathscr{L}^{k+1}$ ), and we set

$$
\delta^{k_{0}-2}=d^{\prime} d^{\prime \prime}: \mathscr{L}^{k_{0}-2}=\mathscr{E}^{p_{0}-1, q_{0}-1} \longrightarrow \mathscr{L}^{k_{0}-1}=\mathscr{E}^{p_{0}, q_{0}} .
$$

We find in particular $H_{\mathrm{BC}}^{p_{0}, q_{0}}(X, \mathbb{C})=H^{k_{0}-1}(\mathscr{L} \bullet(X))$. We observe that $\mathscr{L} \bullet$ has subcomplexes $\left(\mathscr{S}^{\prime \bullet}, d^{\prime}\right)$ and $\left(\mathscr{S}^{\prime \prime \bullet}, d^{\prime \prime}\right)$ defined by

$$
\begin{array}{llll}
\mathscr{S}^{\prime k}=\frac{\Omega_{X}^{k}}{} & \text { for } \quad 0 \leqslant k \leqslant p_{0}-1, & \mathscr{S}^{\prime k}=0 & \text { otherwise, } \\
\mathscr{S}^{\prime \prime k}=\overline{\Omega_{X}^{k}} & \text { for } 0 \leqslant k \leqslant q_{0}-1, & \mathscr{S}^{\prime \prime k}=0 & \text { otherwise. }
\end{array}
$$

If $p_{0}=0$ or $q_{0}=0$ we set instead $\mathscr{S}^{\prime 0}=\mathbb{C}$ or $\mathscr{S}^{\prime \prime 0}=\mathbb{C}$, and take the other components to be zero. Finally, we let $\mathscr{S}^{\bullet}=\mathscr{S}^{\bullet \bullet}+\mathscr{S}^{\prime \prime} \subset \mathscr{L}^{\bullet}$ (the sum is direct except for $\mathscr{S}^{0}$ ); we denote by $\mathscr{M}^{\bullet}$ the sheaf complex defined in the same way as $\mathscr{L}^{\bullet}$, except that the sheaves $\mathscr{E}^{p, q}$ are replaced by the sheaves of currents $\mathscr{D}_{n-p, n-q}^{\prime}$.
(12.1) Lemma. The inclusions $\mathscr{S}^{\bullet} \subset \mathscr{L}^{\bullet} \subset \mathscr{M}^{\bullet}$ induce isomorphisms

$$
\mathscr{H}^{k}\left(\mathscr{S}^{\bullet}\right) \simeq \mathscr{H}^{k}\left(\mathscr{L}^{\bullet}\right) \simeq \mathscr{H}^{k}\left(\mathscr{M}^{\bullet}\right)
$$

and these cohomology sheaves vanish for $k \neq 0, p_{0}-1, q_{0}-1$.
Proof. We will prove the result only for the inclusion $\mathscr{S}^{\bullet} \subset \mathscr{L}^{\bullet}$, the other case $\mathscr{S}^{\bullet} \subset \mathscr{M}^{\bullet}$ is identical. Let us denote by $\mathscr{Z}^{p, q}$ the sheaf of $d^{\prime \prime}$-closed differential forms of bidegree $(p, q)$. We consider the filtration

$$
F_{p}\left(\mathscr{L}^{k}\right)=\mathscr{L}^{k} \cap \bigoplus_{r \geqslant p} \mathscr{E}^{r, \bullet}
$$

and the induced filtration on $\mathscr{S}^{\bullet}$. In the case of $\mathscr{L}^{\bullet}$, the first spectral sequence has the following terms $E_{0}^{\bullet}$ and $E_{1}^{\bullet}$ :

$$
\begin{aligned}
& \text { if } p<p_{0} \quad E_{0}^{p, \bullet}: \quad 0 \longrightarrow \mathscr{E}^{p, 0} \xrightarrow{d^{\prime \prime}} \mathscr{E}^{p, 1} \longrightarrow \cdots \xrightarrow{d^{\prime \prime}} \mathscr{E}^{p, q_{0}-1} \longrightarrow 0, \\
& \text { if } p \geqslant p_{0} \quad E_{0}^{p, \bullet}: 0 \longrightarrow \mathscr{E}^{p, q_{0}} \xrightarrow{d^{\prime \prime}} \mathscr{E}^{p, q_{0}+1} \longrightarrow \cdots \longrightarrow \mathscr{E}^{p, q} \xrightarrow{d^{\prime \prime}} \cdots, \\
& \text { if } p<p_{0} \\
& \text { if } p \geqslant E_{0}^{p, 0}=\Omega_{X}^{p}, E_{1}^{p, q_{0}-1} \simeq E_{1}^{p, q_{0}-1}=\mathscr{E}^{p, q_{0}}, \quad E_{1}^{p, q_{0}}, E_{1}^{p, q}=0 \text { for } q \neq q_{0}-1 .
\end{aligned}
$$

The isomorphism in the third line is given by

$$
\mathscr{E}^{p, q_{0}-1} / d^{\prime \prime} \mathscr{E}^{p, q_{0}-2} \simeq d^{\prime \prime} \mathscr{E}^{p, q_{0}-1} \simeq \mathscr{Z}^{p, q_{0}} .
$$

The map $d_{1}: E_{1}^{p_{0}-1, q_{0}-1} \longrightarrow E_{1}^{p_{0}, q_{0}-1}$ is induced by $d^{\prime} d^{\prime \prime}$ acting on $\mathscr{E}^{p_{0}-1, q_{0}-1}$, but thanks to the previous identification, this map becomes $d^{\prime}$ acting on $\mathscr{L}^{p_{0}-1, q_{0}}$. Hence $E_{1}^{\bullet}$ consists of two sequences

$$
\begin{aligned}
& E_{1}^{\bullet, 0}: 0 \longrightarrow \Omega_{X}^{0} \xrightarrow{d^{\prime}} \Omega_{X}^{1} \longrightarrow \cdots \xrightarrow{d^{\prime}} \Omega_{X}^{p_{0}-1} \longrightarrow 0, \\
& E_{1}^{\bullet, q_{0}-1}: 0 \longrightarrow \mathscr{L}^{0, q_{0}} \xrightarrow{d^{\prime}} \mathscr{L}^{1, q_{0}} \longrightarrow \cdots \longrightarrow \not \mathscr{L}^{p, q_{0}} \xrightarrow{d^{\prime}} \cdots ;
\end{aligned}
$$

if these sequences overlap ( $q_{0}=1$ ), only the second one has to be considered. The term $E_{1}^{\bullet}$ in the spectral sequence of $\mathscr{S}^{\bullet}$ has the same first line, but the second is reduced to $E_{1}^{0, q_{0}-1}=\overline{d \Omega_{X}^{q_{0}-2}}$ (resp. $=\mathbb{C}$ for $q_{0}=1$ ). Thanks to Lemma 12.2 below, we see that the two spectral sequences coincide in $E_{2}^{\bullet}$, with at most three non zero terms:

$$
E_{2}^{0,0}=\mathbb{C}, \quad E_{2}^{p_{0}-1,0}=d \Omega_{X}^{p_{0}-2} \text { for } p_{0} \geqslant 2, \quad E_{2}^{0, q_{0}-1}=\overline{d \Omega_{X}^{q_{0}-2}} \text { for } q_{0} \geqslant 2
$$

Hence $\mathscr{H}^{k}\left(\mathscr{S}^{\bullet}\right) \simeq \mathscr{H}^{k}\left(\mathscr{L}^{\bullet}\right)$ and these sheaves vanish for $k \neq 0, p_{0}-1, q_{0}-1$.
(12.2) Lemma. The complex of sheaves

$$
0 \longrightarrow \mathscr{Z}^{0, q_{0}} \xrightarrow{d^{\prime}} \mathscr{Z}^{1, q_{0}} \longrightarrow \cdots \longrightarrow \mathscr{Z}^{p, q_{0}} \xrightarrow{d^{\prime}} \cdots
$$

is a resolution of $\overline{d \Omega_{X}^{q_{0}-1}}$ for $q_{0} \geqslant 1$, resp. of $\mathbb{C}$ for $q_{0}=0$.
Proof. Embed $\mathscr{L}^{\bullet}, q_{0}$ in the double complex

$$
K^{p, q}=\mathscr{E}^{p, q} \quad \text { for } \quad q<q_{0}, \quad K^{p, q}=0 \quad \text { for } \quad q \geqslant q_{0} .
$$

For the first fitration of $K^{\bullet}$, we find

$$
E_{1}^{p, q_{0}-1}=\mathscr{Z}^{p, q_{0}}, \quad E_{1}^{p, q}=0 \quad \text { for } \quad q \neq q_{0}-1
$$

The second fitration gives $\widetilde{E}_{1}^{p, q}=0$ for $q \geqslant 1$ and

$$
\widetilde{E}_{1}^{p, 0}=H^{0}\left(K^{\bullet, p}\right)= \begin{cases}\overline{H^{0}\left(\mathscr{E}^{p, \bullet}\right)} & \overline{\Omega_{X}^{p}} \\ 0 & \text { for } p \leqslant q_{0}-1 \\ \text { for } p \geqslant q_{0}\end{cases}
$$

thus the cohomology of $\mathscr{\mathscr { E }} \bullet, q_{0}$ coincides with that of $\left(\overline{\Omega_{X}^{p}}, d\right)_{0 \leqslant p<q_{0}}$.

Lemma IV-11.10 and formula (IV-12.9) imply

$$
\begin{align*}
\mathbb{H}^{k}\left(X, \mathscr{S}^{\bullet}\right) & \simeq \mathbb{H}^{k}\left(X, \mathscr{L}^{\bullet}\right) \simeq \mathbb{H}^{k}\left(X, \mathscr{M}^{\bullet}\right)  \tag{12.3}\\
& \simeq H^{k}\left(\mathscr{L}^{\bullet}(X)\right) \simeq H^{k}\left(\mathscr{M}^{\bullet}(X)\right)
\end{align*}
$$

because the sheaves $\mathscr{L}^{k}$ and $\mathscr{M}^{k}$ are soft. In particular, the group $H_{\mathrm{BC}}^{p, q}(X, \mathbb{C})$ can be computed either by means of $\mathscr{C}^{\infty}$ differential forms or by means of currents. This property also holds for the De Rham or Dolbeault groups $H^{k}(X, \mathbb{C}), H^{p, q}(X, \mathbb{C})$, as was already remarked in §IV-6. Another important consequence of (12.3) is:
(12.4) Theorem. If $X$ is compact, then $\operatorname{dim} H_{\mathrm{BC}}^{p, q}(X, \mathbb{C})<+\infty$.

Proof. We show more generally that the hypercohomology groups $\mathbb{H}^{k}\left(X, \mathscr{S}^{\bullet}\right)$ are finite dimensional. As there is an exact sequence

$$
0 \longrightarrow \mathbb{C} \longrightarrow \mathscr{S}^{\prime} \bullet \oplus \mathscr{S}^{\prime \prime} \bullet \longrightarrow \mathscr{S}^{\bullet} \longrightarrow 0
$$

and a corresponding long exact sequence for hypercohomology groups, it is enough to show that the groups $\mathbb{H}^{k}\left(X, \mathscr{S}^{\prime \bullet}\right)$ are finite dimensional. This property is proved for $\mathscr{S}^{\prime \bullet}=\mathscr{S}_{p_{0}}^{\prime \bullet}$ by induction on $p_{0}$. For $p_{0}=0$ or 1 , the complex $\mathscr{S}^{\prime \bullet}$ is reduced to its term $\mathscr{G}^{\prime 0}$, thus

$$
\mathbb{H}^{k}\left(X, \mathscr{S}^{\bullet}\right)=H^{k}\left(X, \mathscr{S}^{\prime 0}\right)= \begin{cases}H^{k}(X, \mathbb{C}) & \text { for } p_{0}=0 \\ H^{k}(X, \mathscr{C}) & \text { for } p_{0}=1\end{cases}
$$

and this groups are finite dimensional. In general, we have an exact sequence

$$
0 \longrightarrow \Omega_{X}^{p_{0}} \longrightarrow \mathscr{S}_{p_{0}+1}^{\bullet} \longrightarrow \mathscr{S}_{p_{0}}^{\bullet} \longrightarrow 0
$$

where $\Omega_{X}^{p_{0}}$ denotes the subcomplex of $\mathscr{S}_{p_{0}+1}^{\bullet}$ reduced to one term in degree $p_{0}$. As

$$
\mathbb{H}^{k}\left(X, \Omega_{X}^{p_{0}}\right)=H^{k-p_{0}}\left(X, \Omega_{X}^{p_{0}}\right)=H^{p_{0}, k-p_{0}}(X, \mathbb{C})
$$

is finite dimensional, the Theorem follows.
(12.5) Definition. We say that a compact manifold admits a strong Hodge decomposition if the natural maps

$$
H_{\mathrm{BC}}^{p, q}(X, \mathbb{C}) \longrightarrow H^{p, q}(X, \mathbb{C}), \quad \bigoplus_{p+q=k} H_{\mathrm{BC}}^{p, q}(X, \mathbb{C}) \longrightarrow H^{k}(X, \mathbb{C})
$$

are isomorphisms.

This implies of course that there are natural isomorphisms

$$
H^{k}(X, \mathbb{C}) \simeq \bigoplus_{p+q=k} H^{p, q}(X, \mathbb{C}), \quad H^{q, p}(X, \mathbb{C}) \simeq \overline{H^{p, q}(X, \mathbb{C})}
$$

and that the Hodge-Frölicher spectral sequence degenerates in $E_{1}^{\bullet}$. It follows from $\S 8$ that all Kähler manifolds admit a strong Hodge decomposition.

## § 12.2. Direct and Inverse Image Morphisms

Let $F: X \longrightarrow Y$ be a holomorphic map between complex analytic manifolds of respective dimensions $n, m$, and $r=n-m$. We have pull-back morphisms

$$
\begin{align*}
& F^{\star}: \quad H^{k}(Y, \mathbb{C}) \longrightarrow H^{k}(X, \mathbb{C}), \\
& F^{\star}: \quad H^{p, q}(Y, \mathbb{C}) \longrightarrow H^{p, q}(X, \mathbb{C}),  \tag{12.6}\\
& F^{\star}: H_{\mathrm{BC}}^{p, q}(Y, \mathbb{C}) \longrightarrow H_{\mathrm{BC}}^{p, q}(X, \mathbb{C}),
\end{align*}
$$

commuting with the natural morphisms (8.2), (8.3).
Assume now that $F$ is proper. Theorem I-1.14 shows that one can define direct image morphisms

$$
F_{\star}: \mathscr{D}_{k}^{\prime}(X) \longrightarrow \mathscr{D}_{k}^{\prime}(Y), \quad F_{\star}: \quad \mathscr{D}_{p, q}^{\prime}(X) \longrightarrow \mathscr{D}_{p, q}^{\prime}(Y),
$$

commuting with $d^{\prime}, d^{\prime \prime}$. To $F_{\star}$ therefore correspond cohomology morphisms

$$
\begin{align*}
& F_{\star}: \quad H^{k}(X, \mathbb{C}) \longrightarrow H^{k-2 r}(Y, \mathbb{C}), \\
& F_{\star}: \quad H^{p, q}(X, \mathbb{C}) \longrightarrow H^{p-r, q-r}(Y, \mathbb{C}),  \tag{12.7}\\
& F_{\star}: H_{\mathrm{BC}}^{p, q}(X, \mathbb{C}) \longrightarrow H_{\mathrm{BC}}^{p-r, q-r}(Y, \mathbb{C}),
\end{align*}
$$

which commute also with (8.2), (8.3). In addition, I-1.14 c) implies the adjunction formula

$$
\begin{equation*}
F_{\star}\left(\alpha \smile F^{\star} \beta\right)=\left(F_{\star} \alpha\right) \smile \beta \tag{12.8}
\end{equation*}
$$

whenever $\alpha$ is a cohomology class (of any of the three above types) on $X$, and $\beta$ a cohomology class (of the same type) on $Y$.

## § 12.3. Modifications and the Fujiki Class (C)

Recall that a modification of a compact manifold $X$ is a holomorphic map $\mu: \widetilde{X} \longrightarrow X$ such that
i) $\widetilde{X}$ is a compact complex manifold of the same dimension as $X$;
ii) there exists an analytic subset $S \subset X$ of codimension $\geqslant 1$ such that $\mu: \widetilde{X} \backslash \mu^{-1}(S) \longrightarrow$ $X \backslash S$ is a biholomorphism.
(12.9) Theorem. If $\widetilde{X}$ admits a strong Hodge decomposition, and if $\mu: \widetilde{X} \longrightarrow X$ is a modification, then $X$ also admits a strong Hodge decomposition.

Proof. We first observe that $\mu_{\star} \mu^{\star} f=f$ for every smooth form $f$ on $Y$. In fact, this property is equivalent to the equality

$$
\int_{Y}\left(\mu_{\star} \mu^{\star} f\right) \wedge g=\int_{X} \mu^{\star}(f \wedge g)=\int_{Y} f \wedge g
$$

for every smooth form $g$ on $Y$, and this equality is clear because $\mu$ is a biholomorphism outside sets of Lebesgue measure 0 . Consequently, the induced cohomology morphism $\mu_{\star}$ is surjective and $\mu^{\star}$ is injective (but these maps need not be isomorphisms). Now, we have commutative diagrams

$$
\begin{array}{ccc}
H_{\mathrm{BC}}^{p, q}(\widetilde{X}, \mathbb{C}) \longrightarrow H^{p, q}(\widetilde{X}, \mathbb{C}), & \bigoplus_{p+q=k} H_{\mathrm{BC}}^{p, q}(\widetilde{X}, \mathbb{C}) \longrightarrow H^{k}(\widetilde{X}, \mathbb{C}) \\
\mu_{\star} \downarrow \uparrow \mu^{\star} & \mu_{\star} \downarrow \uparrow \mu^{\star} & \mu_{\star} \downarrow \uparrow \mu^{\star}
\end{array} \mu_{\star} \downarrow \uparrow \mu^{\star} .
$$

with either upward or downward vertical arrows. Hence the surjectivity or injectivity of the top horizontal arrows implies that of the bottom horizontal arrows.
(12.10) Definition. A manifold $X$ is said to be in the Fujiki class (C) if $X$ admits a Kähler modification $\widetilde{X}$.

By Th. 12.9, Hodge decomposition still holds for a manifold in the class (C6). We will see later that there exist non-Kähler manifolds in ( $\mathscr{C}$ ), for example all non projective Moišezon manifolds (cf. §?.?). The class (C) has been first introduced in [Fujiki 1978].

## Chapter VII

## Positive Vector Bundles and Vanishing Theorems


#### Abstract

In this chapter, we prove a few vanishing theorems for hermitian vector bundles over compact complex manifolds. All these theorems are based on an a priori inequality for $(p, q)$-forms with values in a vector bundle, known as the Bochner-Kodaira-Nakano inequality. This inequality naturally leads to several positivity notions for the curvature of a vector bundle ([Kodaira 1953, 1954], [Griffiths 1969] and [Nakano 1955, 1973]). The corresponding algebraic notion of ampleness introduced by Grothendieck [EGA 196067] and developped further by [Nakai 1963] and [Hartshorne 1966] is also discussed. The differential geometric techniques yield optimal vanishing results in the case of line bundles (Kodaira-Akizuki-Nakano and Girbau vanishing theorems) and also some partial results in the case of vector bundles (Nakano vanishing theorem). As an illustration, we compute the cohomology groups $H^{p, q}\left(\mathbb{P}^{n}, \mathfrak{O}(k)\right)$; much finer results will be obtained in chapters $8-11$. Finally, the Kodaira vanishing theorem is combined with a blowing-up technique in order to establish the projective embedding theorem for manifolds admitting a Hodge metric.


## § 1. Bochner-Kodaira-Nakano Identity

Let $(X, \omega)$ be a hermitian manifold, $\operatorname{dim}_{\mathbb{C}} X=n$, and let $E$ be a hermitian holomorphic vector bundle of rank $r$ over $X$. We denote by $D=D^{\prime}+D^{\prime \prime}$ its Chern connection (or $D_{E}$ if we want to specify the bundle), and by $\delta=\delta^{\prime}+\delta^{\prime \prime}$ the formal adjoint operator of $D$. The operators $L, \Lambda$ of chapter 6 are extended to vector valued forms in $\Lambda^{p, q} T^{\star} X \otimes E$ by taking their tensor product with $\operatorname{Id}_{E}$. The following result extends the commutation relations of chapter 6 to the case of bundle valued operators.
(1.1) Theorem. If $\tau$ is the operator of type $(1,0)$ defined by $\tau=\left[\Lambda, d^{\prime} \omega\right]$ on $\mathscr{C}_{\bullet, \bullet}^{\infty}(X, E)$, then
a) $\quad\left[\delta_{E}^{\prime \prime}, L\right]=\mathrm{i}\left(D_{E}^{\prime}+\tau\right)$,
b) $\left[\delta_{E}^{\prime}, L\right]=-\mathrm{i}\left(D_{E}^{\prime \prime}+\bar{\tau}\right)$,
c) $\left[\Lambda, D_{E}^{\prime \prime}\right]=-\mathrm{i}\left(\delta_{E}^{\prime}+\tau^{\star}\right)$,
d) $\left[\Lambda, D_{E}^{\prime}\right]=\mathrm{i}\left(\delta_{E}^{\prime \prime}+\bar{\tau}^{\star}\right)$.

Proof. Fix a point $x_{0}$ in X and a coordinate system $z=\left(z_{1}, \ldots, z_{n}\right)$ centered at $x_{0}$. Then Prop. V-12.?? shows the existence of a normal coordinate frame ( $e_{\lambda}$ ) at $x_{0}$. Given any section $s=\sum_{\lambda} \sigma_{\lambda} \otimes e_{\lambda} \in \mathscr{C}_{p, q}^{\infty}(X, E)$, it is easy to check that the operators $D_{E}, \delta_{E}^{\prime \prime}, \ldots$ have Taylor expansions of the type

$$
D_{E} s=\sum_{\lambda} d \sigma_{\lambda} \otimes e_{\lambda}+O(|z|), \quad \delta_{E}^{\prime \prime} s=\sum_{\lambda} \delta^{\prime \prime} \sigma_{\lambda} \otimes e_{\lambda}+O(|z|), \ldots
$$

in terms of the scalar valued operators $d, \delta, \ldots$ Here the terms $O(|z|)$ depend on the curvature coefficients of $E$. The proof of Th. 1.1 is then reduced to the case of scalar valued operators, which is granted by Th. VI-10.1.

The Bochner-Kodaira-Nakano identity expresses the antiholomorphic Laplace operator $\Delta^{\prime \prime}=D^{\prime \prime} \delta^{\prime \prime}+\delta^{\prime \prime} D^{\prime \prime}$ acting on $\mathscr{C}_{\bullet \bullet \bullet}^{\infty}(X, E)$ in terms of its conjugate operator $\Delta^{\prime}=D^{\prime} \delta^{\prime}+\delta^{\prime} D^{\prime}$, plus some extra terms involving the curvature of $E$ and the torsion of the metric $\omega$ (in case $\omega$ is not Kähler). Such identities appear frequently in riemannian geometry (Weitzenböck formula).
(1.2) Theorem. $\Delta^{\prime \prime}=\Delta^{\prime}+[\mathrm{i} \Theta(E), \Lambda]+\left[D^{\prime}, \tau^{\star}\right]-\left[D^{\prime \prime}, \bar{\tau}^{\star}\right]$.

Proof. Equality 1.1 d) yields $\delta^{\prime \prime}=-\mathrm{i}\left[\Lambda, D^{\prime}\right]-\bar{\tau}^{\star}$, hence

$$
\Delta^{\prime \prime}=\left[D^{\prime \prime}, \delta^{\prime \prime}\right]=-\mathrm{i}\left[D^{\prime \prime},\left[\Lambda, D^{\prime}\right]\right]-\left[D^{\prime \prime}, \bar{\tau}^{\star}\right] .
$$

The Jacobi identity VI-10.2 and relation 1.1 c) imply

$$
\left[D^{\prime \prime},\left[\Lambda, D^{\prime}\right]\right]=\left[\Lambda,\left[D^{\prime}, D^{\prime \prime}\right]\right]+\left[D^{\prime},\left[D^{\prime \prime}, \Lambda\right]\right]=[\Lambda, \Theta(E)]+\mathrm{i}\left[D^{\prime}, \delta^{\prime}+\tau^{\star}\right]
$$

taking into account that $\left[D^{\prime}, D^{\prime \prime}\right]=D^{2}=\Theta(E)$. Theorem 1.2 follows.
(1.3) Corollary ([Akizuki-Nakano 1955]). If $\omega$ is Kähler, then

$$
\Delta^{\prime \prime}=\Delta^{\prime}+[\mathrm{i} \Theta(E), \Lambda] .
$$

In the latter case, $\Delta^{\prime \prime}-\Delta^{\prime}$ is therefore an operator of order 0 closely related to the curvature of $E$. When $\omega$ is not Kähler, Formula 1.2 is not really satisfactory, because it involves the first order operators $\left[D^{\prime}, \tau^{\star}\right]$ and $\left[D^{\prime \prime}, \bar{\tau}^{\star}\right]$. In fact, these operators can be combined with $\Delta^{\prime}$ in order to yield a new positive self-adjoint operator $\Delta_{\tau}^{\prime}$.
(1.4) Theorem ([Demailly 1985]). The operator $\Delta_{\tau}^{\prime}=\left[D^{\prime}+\tau, \delta^{\prime}+\tau^{\star}\right]$ is a positive and formally self-adjoint operator with the same principal part as the Laplace operator $\Delta^{\prime}$. Moreover

$$
\Delta^{\prime \prime}=\Delta_{\tau}^{\prime}+[\mathrm{i} \Theta(E), \Lambda]+T_{\omega}
$$

where $T_{\omega}$ is an operator of order 0 depending only on the torsion of the hermitian metric $\omega$ :

$$
T_{\omega}=\left[\Lambda,\left[\Lambda, \frac{\mathrm{i}}{2} d^{\prime} d^{\prime \prime} \omega\right]\right]-\left[d^{\prime} \omega,\left(d^{\prime} \omega\right)^{\star}\right] .
$$

Proof. The first assertion is clear, because the equality $\left(D^{\prime}+\tau\right)^{\star}=\delta^{\prime}+\tau^{\star}$ implies the self-adjointness of $\Delta_{\tau}^{\prime}$ and

$$
\left\langle\left\langle\Delta_{\tau}^{\prime} u, u\right\rangle\right\rangle=\left\|D^{\prime} u+\tau u\right\|^{2}+\left\|\delta^{\prime} u+\tau^{\star} u\right\|^{2} \geqslant 0
$$

for any compactly supported form $u \in \mathscr{C}_{p, q}^{\infty}(X, E)$. In order to prove the formula, we need two lemmas.
(1.5) Lemma. a)

$$
[L, \tau]=3 d^{\prime} \omega
$$

b) $\quad[\Lambda, \tau]=-2 \mathrm{i} \bar{\tau}^{\star}$.

Proof. a) Since $\left[L, d^{\prime} \omega\right]=0$, the Jacobi identity implies

$$
[L, \tau]=\left[L,\left[\Lambda, d^{\prime} \omega\right]\right]=-\left[d^{\prime} \omega,[L, \Lambda]\right]=3 d^{\prime} \omega
$$

taking into account Cor. VI-10.4 and the fact that $d^{\prime} \omega$ is of degree 3 .
b) By 1.1 a) we have $\tau=-\mathrm{i}\left[\delta^{\prime \prime}, L\right]-D^{\prime}$, hence

$$
[\Lambda, \tau]=-i\left[\Lambda,\left[\delta^{\prime \prime}, L\right]\right]-\left[\Lambda, D^{\prime}\right]=-i\left(\left[\Lambda,\left[\delta^{\prime \prime}, L\right]\right]+\delta^{\prime \prime}+\bar{\tau}^{\star}\right)
$$

Using again VI-10.4 and the Jacobi identity, we get

$$
\begin{aligned}
{\left[\Lambda,\left[\delta^{\prime \prime}, L\right]\right] } & =-\left[L,\left[\Lambda, \delta^{\prime \prime}\right]\right]-\left[\delta^{\prime \prime},[L, \Lambda]\right] \\
& =-\left[\left[d^{\prime \prime}, L\right], \Lambda\right]^{\star}-\delta^{\prime \prime}=-\left[d^{\prime \prime} \omega, \Lambda\right]^{\star}-\delta^{\prime \prime}=\bar{\tau}^{\star}-\delta^{\prime \prime}
\end{aligned}
$$

A substitution in the previous equality gives $[\Lambda, \tau]=-2 \mathrm{i} \bar{\tau}^{\star}$.
(1.6) Lemma. The following identities hold:
a) $\left[D^{\prime}, \bar{\tau}^{\star}\right]=-\left[D^{\prime}, \delta^{\prime \prime}\right]=\left[\tau, \delta^{\prime \prime}\right]$,
b) $-\left[D^{\prime \prime}, \bar{\tau}^{\star}\right]=\left[\tau, \delta^{\prime}+\tau^{\star}\right]+T_{\omega}$.

Proof. a) The Jacobi identity implies

$$
-\left[D^{\prime},\left[\Lambda, D^{\prime}\right]\right]+\left[D^{\prime},\left[D^{\prime}, \Lambda\right]\right]+\left[\Lambda,\left[D^{\prime}, D^{\prime}\right]\right]=0
$$

hence $-2\left[D^{\prime},\left[\Lambda, D^{\prime}\right]\right]=0$ and likewise $\left[\delta^{\prime \prime},\left[\delta^{\prime \prime}, L\right]\right]=0$. Assertion a) is now a consequence of 1.1 a ) and d$)$.
b) In order to verify b), we start from the equality $\bar{\tau}^{\star}=\frac{i}{2}[\Lambda, \tau]$ provided by Lemma 1.5 b). It follows that

$$
\begin{equation*}
\left[D^{\prime \prime}, \bar{\tau}^{\star}\right]=\frac{\mathrm{i}}{2}\left[D^{\prime \prime},[\Lambda, \tau]\right] \tag{1.7}
\end{equation*}
$$

The Jacobi identity will now be used several times. One obtains

$$
\begin{align*}
& {\left[D^{\prime \prime},[\Lambda, \tau]\right]=\left[\Lambda,\left[\tau, D^{\prime \prime}\right]\right]+\left[\tau,\left[D^{\prime \prime}, \Lambda\right]\right]}  \tag{1.8}\\
& {\left[\tau, D^{\prime \prime}\right]=\left[D^{\prime \prime}, \tau\right]=\left[D^{\prime \prime},\left[\Lambda, d^{\prime} \omega\right]\right]=\left[\Lambda,\left[d^{\prime} \omega, D^{\prime \prime}\right]\right]+\left[d^{\prime} \omega,\left[D^{\prime \prime}, \Lambda\right]\right]}  \tag{1.9}\\
& \quad=\left[\Lambda, d^{\prime \prime} d^{\prime} \omega\right]+\left[d^{\prime} \omega, A\right]
\end{align*}
$$

with $A=\left[D^{\prime \prime}, \Lambda\right]=\mathrm{i}\left(\delta^{\prime}+\tau^{\star}\right)$. From (1.9) we deduce

$$
\begin{equation*}
\left[\Lambda,\left[\tau, D^{\prime \prime}\right]\right]=\left[\Lambda,\left[\Lambda, d^{\prime \prime} d^{\prime} \omega\right]\right]+\left[\Lambda,\left[d^{\prime} \omega, A\right]\right] \tag{1.10}
\end{equation*}
$$

Let us compute now the second Lie bracket in the right hand side of (1.10:

$$
\begin{align*}
{\left[\Lambda,\left[d^{\prime} \omega, A\right]\right] } & =\left[A,\left[\Lambda, d^{\prime} \omega\right]\right]-\left[d^{\prime} \omega,[A, \Lambda]\right]=[\tau, A]+\left[d^{\prime} \omega,[\Lambda, A]\right] ;  \tag{1.11}\\
{[\Lambda, A] } & =\mathrm{i}\left[\Lambda, \delta^{\prime}+\tau^{\star}\right]=\mathrm{i}\left[D^{\prime}+\tau, L\right]^{\star} . \tag{1.12}
\end{align*}
$$

Lemma 1.5 b ) provides $[\tau, L]=-3 d^{\prime} \omega$, and it is clear that $\left[D^{\prime}, L\right]=d^{\prime} \omega$. Equalities (1.12) and (1.11) yield therefore

$$
\begin{align*}
{[\Lambda, A] } & =-2 \mathrm{i}\left(d^{\prime} \omega\right)^{\star} \\
{\left[\Lambda,\left[d^{\prime} \omega, A\right]\right] } & =\left[\tau,\left[D^{\prime \prime}, \Lambda\right]\right]-2 \mathrm{i}\left[d^{\prime} \omega,\left(d^{\prime} \omega\right)^{\star}\right] . \tag{1.13}
\end{align*}
$$

Substituting (1.10) and (1.13) in (1.8) we get

$$
\begin{align*}
{\left[D^{\prime \prime},[\Lambda, \tau]\right] } & =\left[\Lambda,\left[\Lambda, d^{\prime \prime} d^{\prime} \omega\right]\right]+2\left[\tau,\left[D^{\prime \prime}, \Lambda\right]\right]-2 \mathrm{i}\left[d^{\prime} \omega,\left(d^{\prime} \omega\right)^{\star}\right]  \tag{1.14}\\
& =2 \mathrm{i}\left(T_{\omega}+\left[\tau, \delta^{\prime}+\tau^{\star}\right]\right)
\end{align*}
$$

Formula b) is a consequence of (1.7) and (1.14).
Theorem 1.4 follows now from Th. 1.2 if Formula 1.6 b) is rewritten

$$
\Delta^{\prime}+\left[D^{\prime}, \tau^{\star}\right]-\left[D^{\prime \prime}, \bar{\tau}^{\star}\right]=\left[D^{\prime}+\tau, \delta^{\prime}+\tau^{\star}\right]+T_{\omega}
$$

When $\omega$ is Kähler, then $\tau=T_{\omega}=0$ and Lemma 1.6 a) shows that $\left[D^{\prime}, \delta^{\prime \prime}\right]=0$. Together with the adjoint relation $\left[D^{\prime \prime}, \delta^{\prime}\right]=0$, this equality implies

$$
\begin{equation*}
\Delta=\Delta^{\prime}+\Delta^{\prime \prime} \tag{1.15}
\end{equation*}
$$

When $\omega$ is not Kähler, Lemma 1.6 a) can be written $\left[D^{\prime}+\tau, \delta^{\prime \prime}\right]=0$ and we obtain more generally

$$
\left[D+\tau, \delta+\tau^{\star}\right]=\left[\left(D^{\prime}+\tau\right)+D^{\prime \prime},\left(\delta^{\prime}+\tau^{\star}\right)+\delta^{\prime \prime}\right]=\Delta_{\tau}^{\prime}+\Delta^{\prime \prime}
$$

(1.16) Proposition. Set $\Delta_{\tau}=\left[D+\tau, \delta+\tau^{\star}\right]$. Then $\Delta_{\tau}=\Delta_{\tau}^{\prime}+\Delta^{\prime \prime}$.

## § 2. Basic a Priori Inequality

Let $(X, \omega)$ be a compact hermitian manifold, $\operatorname{dim}_{\mathbb{C}} X=n$, and $E$ a hermitian holomorphic vector bundle over $X$. For any section $u \in \mathscr{C}_{p, q}^{\infty}(X, E)$ we have $\left\langle\left\langle\Delta^{\prime \prime} u, u\right\rangle\right\rangle=$ $\left\|D^{\prime \prime} u\right\|^{2}+\left\|\delta^{\prime \prime} u\right\|^{2}$ and the similar formula for $\Delta_{\tau}^{\prime}$ gives $\left\langle\left\langle\Delta_{\tau}^{\prime} u, u\right\rangle\right\rangle \geqslant 0$. Theorem 1.4 implies therefore

$$
\begin{equation*}
\left\|D^{\prime \prime} u\right\|^{2}+\left\|\delta^{\prime \prime} u\right\|^{2} \geqslant \int_{X}\left(\langle[\mathrm{i} \Theta(E), \Lambda] u, u\rangle+\left\langle T_{\omega} u, u\right\rangle\right) d V . \tag{2.1}
\end{equation*}
$$

This inequality is known as the Bochner-Kodaira-Nakano inequality. When $u$ is $\Delta^{\prime \prime}$ harmonic, we get in particular

$$
\begin{equation*}
\int_{X}\left(\langle[\mathrm{i} \Theta(E), \Lambda] u, u\rangle+\left\langle T_{\omega} u, u\right\rangle\right) d V \leqslant 0 \tag{2.2}
\end{equation*}
$$

These basic a priori estimates are the starting point of all vanishing theorems. Observe that $[\mathrm{i} \Theta(E), \Lambda]+T_{\omega}$ is a hermitian operator acting pointwise on $\Lambda^{p, q} T^{\star} X \otimes E$ (the
hermitian property can be seen from the fact that this operator coincides with $\Delta^{\prime \prime}-\Delta_{\tau}^{\prime}$ on smooth sections). Using Hodge theory (Cor. VI-11.2), we get:
(2.3) Corollary. If the hermitian operator $[\mathrm{i} \Theta(E), \Lambda]+T_{\omega}$ is positive definite on $\Lambda^{p, q} T^{\star} X \otimes E$, then $H^{p, q}(X, E)=0$.

In some circumstances, one can improve Cor. 2.3 thanks to the following "analytic continuation lemma" due to [Aronszajn 1957]:
(2.4) Lemma. Let $M$ be a connected $\mathscr{C}^{\infty}$-manifold, $F$ a vector bundle over $M$, and $P$ a second order elliptic differential operator acting on $\mathscr{C}^{\infty}(M, F)$. Then any section $\alpha \in \operatorname{ker} P$ vanishing on a non-empty open subset of $M$ vanishes identically on $M$.
(2.5) Corollary. Assume that $X$ is compact and connected. If

$$
[\mathrm{i} \Theta(E), \Lambda]+T_{\omega} \in \operatorname{Herm}\left(\Lambda^{p, q} T^{\star} X \otimes E\right)
$$

is semi-positive on $X$ and positive definite in at least one point $x_{0} \in X$, then we have $H^{p, q}(X, E)=0$.

Proof. By (2.2) every $\Delta^{\prime \prime}$-harmonic ( $p, q$ )-form $u$ must vanish in the neighborhood of $x_{0}$ where $[\mathrm{i} \Theta(E), \Lambda]+T_{\omega}>0$, thus $u \equiv 0$. Hodge theory implies $H^{p, q}(X, E)=0$.

## § 3. Kodaira-Akizuki-Nakano Vanishing Theorem

The main goal of vanishing theorems is to find natural geometric or algebraic conditions on a bundle $E$ that will ensure that some cohomology groups with values in $E$ vanish. In the next three sections, we prove various vanishing theorems for cohomology groups of a hermitian line bundle $E$ over a compact complex manifold $X$.
(3.1) Definition. A hermitian holomorphic line bundle $E$ on $X$ is said to be positive (resp. negative) if the hermitian matrix $\left(c_{j k}(z)\right)$ of its Chern curvature form

$$
\mathrm{i} \Theta(E)=\mathrm{i} \sum_{1 \leqslant j, k \leqslant n} c_{j k}(z) d z_{j} \wedge d \bar{z}_{k}
$$

is positive (resp. negative) definite at every point $z \in X$.

Assume that $X$ has a Kähler metric $\omega$. Let

$$
\gamma_{1}(x) \leqslant \ldots \leqslant \gamma_{n}(x)
$$

be the eigenvalues of $\mathrm{i} \Theta(E)_{x}$ with respect to $\omega_{x}$ at each point $x \in X$, and let

$$
\mathrm{i} \Theta(E)_{x}=\mathrm{i} \sum_{1 \leqslant j \leqslant n} \gamma_{j}(x) \zeta_{j} \wedge \bar{\zeta}_{j}, \quad \zeta_{j} \in T_{x}^{\star} X
$$

be a diagonalization of $\mathrm{i} \Theta(E)_{x}$. By Prop. VI- 8.3 we have

$$
\begin{align*}
\langle[\mathrm{i} \Theta(E), \Lambda] u, u\rangle & =\sum_{J, K}\left(\sum_{j \in J} \gamma_{j}+\sum_{j \in K} \gamma_{j}-\sum_{1 \leqslant j \leqslant n} \gamma_{j}\right)\left|u_{J, K}\right|^{2} \\
& \geqslant\left(\gamma_{1}+\ldots+\gamma_{q}-\gamma_{p+1}-\ldots-\gamma_{n}\right)|u|^{2} \tag{3.2}
\end{align*}
$$

for any form $u=\sum_{J, K} u_{J, K} \zeta_{J} \wedge \bar{\zeta}_{K} \in \Lambda^{p, q} T^{\star} X$.
(3.3) Akizuki-Nakano vanishing theorem (1954). Let $E$ be a holomorphic line bundle on $X$.
a) If $E$ is positive, then $H^{p, q}(X, E)=0$ for $p+q \geqslant n+1$.
b) If $E$ is negative, then $H^{p, q}(X, E)=0$ for $p+q \leqslant n-1$.

Proof. In case a), choose $\omega=\mathrm{i} \Theta(E)$ as a Kähler metric on $X$. Then we have $\gamma_{j}(x)=1$ for all $j$ and $x$, so that

$$
\langle\langle[\mathrm{i} \Theta(E), \Lambda] u, u\rangle\rangle \geqslant(p+q-n)\|u\|^{2}
$$

for any $u \in \Lambda^{p, q} T^{\star} X \otimes E$. Assertion a) follows now from Corollary 2.3. Property b) is proved similarly, by taking $\omega=-\mathrm{i} \Theta(E)$. One can also derive b) from a) by Serre duality (Theorem VI-11.3).

When $p=0$ or $p=n$, Th. 3.3 can be generalized to the case where $\mathrm{i} \Theta(E)$ degenerates at some points. We use here the standard notations

$$
\begin{equation*}
\Omega_{X}^{p}=\Lambda^{p} T^{\star} X, \quad K_{X}=\Lambda^{n} T^{\star} X, \quad n=\operatorname{dim}_{\mathbb{C}} X \tag{3.4}
\end{equation*}
$$

$K_{X}$ is called the canonical line bundle of $X$.
(3.5) Theorem ([Grauert-Riemenschneider 1970]). Let $(X, \omega)$ be a compact and connected Kähler manifold and $E$ a line bundle on $X$.
a) If $\mathrm{i} \Theta(E) \geqslant 0$ on $X$ and $\mathrm{i} \Theta(E)>0$ in at least one point $x_{0} \in X$, then $H^{q}\left(X, K_{X} \otimes E\right)=0 \quad$ for $\quad q \geqslant 1$.
b) If $\mathrm{i} \Theta(E) \leqslant 0$ on $X$ and $\mathrm{i} \Theta(E)<0$ in at least one point $x_{0} \in X$, then

$$
H^{q}(X, E)=0 \quad \text { for } \quad q \leqslant n-1
$$

It will be proved in Volume II, by means of holomorphic Morse inequalities, that the Kähler assumption is in fact unnecessary. This improvement is a deep result first proved by [Siu 1984] with a different ad hoc method.

Proof. For $p=n$, formula (3.2) gives

$$
\begin{equation*}
\langle\langle[i \Theta(E), \Lambda] u, u\rangle\rangle \geqslant\left(\gamma_{1}+\ldots+\gamma_{q}\right)|u|^{2} \tag{3.6}
\end{equation*}
$$

and a) follows from Cor. 2.5. Now b) is a consequence of a) by Serre duality.

## § 4. Girbau's Vanishing Theorem

Let $E$ be a line bundle over a compact connected Kähler manifold $(X, \omega)$. Girbau's theorem deals with the (possibly everywhere) degenerate semi-positive case. We first state the corresponding generalization of Th. 4.5.
(4.1) Theorem. If $\mathrm{i} \Theta(E)$ is semi-positive and has at least $n-s+1$ positive eigenvalues at a point $x_{0} \in X$ for some integer $s \in\{1, \ldots, n\}$, then

$$
H^{q}\left(X, K_{X} \otimes E\right)=0 \quad \text { for } \quad q \geqslant s
$$

Proof. Apply 2.5 and inequality (3.6), and observe that $\gamma_{q}\left(x_{0}\right)>0$ for all $q \geqslant s$.
(4.2) Theorem ([Girbau 1976]). If $\mathrm{i} \Theta(E)$ is semi-positive and has at least $n-s+1$ positive eigenvalues at every point $x \in X$, then

$$
H^{p, q}(X, E)=0 \quad \text { for } \quad p+q \geqslant n+s
$$

Proof. Let us consider on $X$ the new Kähler metric

$$
\omega_{\varepsilon}=\varepsilon \omega+\mathrm{i} \Theta(E), \quad \varepsilon>0
$$

and let $\mathrm{i} \Theta(E)=\mathrm{i} \sum \gamma_{j} \zeta_{j} \wedge \bar{\zeta}_{j}$ be a diagonalization of $\mathrm{i} \Theta(E)$ with respect to $\omega$ and with $\gamma_{1} \leqslant \ldots \leqslant \gamma_{n}$. Then

$$
\omega_{\varepsilon}=\mathrm{i} \sum\left(\varepsilon+\gamma_{j}\right) \zeta_{j} \wedge \bar{\zeta}_{j}
$$

The eigenvalues of $\mathrm{i} \Theta(E)$ with respect to $\omega_{\varepsilon}$ are given therefore by

$$
\begin{equation*}
\gamma_{j, \varepsilon}=\gamma_{j} /\left(\varepsilon+\gamma_{j}\right) \in[0,1[, \quad 1 \leqslant j \leqslant n \tag{4.3}
\end{equation*}
$$

On the other hand, the hypothesis is equivalent to $\gamma_{s}>0$ on $X$. For $j \geqslant s$ we have $\gamma_{j} \geqslant \gamma_{s}$, thus

$$
\begin{equation*}
\gamma_{j, \varepsilon}=\frac{1}{1+\varepsilon / \gamma_{j}} \geqslant \frac{1}{1+\varepsilon / \gamma_{s}} \geqslant 1-\varepsilon / \gamma_{s}, \quad s \leqslant j \leqslant n . \tag{4.4}
\end{equation*}
$$

Let us denote the operators and inner products associated to $\omega_{\varepsilon}$ with $\varepsilon$ as an index. Then inequality (3.2) combined with (4.4) implies

$$
\begin{aligned}
\left\langle\left[\mathrm{i} \Theta(E), \Lambda_{\varepsilon}\right] u, u\right\rangle_{\varepsilon} & \left.\geqslant((q-s+1))\left(1-\varepsilon / \gamma_{s}\right)-(n-p)\right)|u|^{2} \\
& =\left(p+q-n-s+1-(q-s+1) \varepsilon / \gamma_{s}\right)|u|^{2} .
\end{aligned}
$$

Theorem 4.2 follows now from Cor. 2.3 if we choose

$$
\varepsilon<\frac{p+q-n-s+1}{q-s+1} \min _{x \in X} \gamma_{s}(x) .
$$

(4.5) Remark. The following example due to [Ramanujam 1972, 1974] shows that Girbau's result is no longer true for $p<n$ when $\mathrm{i} \Theta(E)$ is only assumed to have $n-s+1$ positive eigenvalues on a dense open set.

Let $V$ be a hermitian vector space of dimension $n+1$ and $X$ the manifold obtained from $P(V) \simeq \mathbb{P}^{n}$ by blowing-up one point $a$. The manifold $X$ may be described as follows: if $P(V / \mathbb{C} a)$ is the projective space of lines $\ell$ containing $a$, then

$$
X=\{(x, \ell) \in P(V) \times P(V / \mathbb{C} a) ; x \in \ell\}
$$

We have two natural projections

$$
\begin{aligned}
& \pi_{1}: X \longrightarrow P(V) \simeq \mathbb{P}^{n}, \\
& \pi_{2}: X \longrightarrow Y=P(V / \mathbb{C} a) \simeq \mathbb{P}^{n-1}
\end{aligned}
$$

It is clear that the preimage $\pi_{1}^{-1}(x)$ is the single point $(x, \ell=(a x))$ if $x \neq a$ and that $\pi_{1}^{-1}(a)=\{a\} \times Y \simeq \mathbb{P}^{n-1}$, therefore

$$
\pi_{1}: X \backslash(\{a\} \times Y) \longrightarrow P(V) \backslash\{a\}
$$

is an isomorphism. On the other hand, $\pi_{2}$ is a locally trivial fiber bundle over $Y$ with fiber $\pi_{2}^{-1}(\ell)=\ell \simeq \mathbb{P}^{1}$, in particular $X$ is smooth and $n$-dimensional. Consider now the line bundle $E=\pi_{1}^{\star} \mathscr{O}(1)$ over $X$, with the hermitian metric induced by that of $\mathbb{O}(1)$. Then $E$ is semi-positive and $\mathrm{i} \Theta(E)$ has $n$ positive eigenvalues at every point of $X \backslash(\{a\} \times Y)$, hence the assumption of Th. 4.2 is satisfied on $X \backslash(\{a\} \times Y)$. However, we will see that

$$
H^{p, p}(X, E) \neq 0, \quad 0 \leqslant p \leqslant n-1,
$$

in contradiction with the expected generalization of (4.2) when $2 p \geqslant n+1$. Let $j: Y \simeq$ $\{a\} \times Y \longrightarrow X$ be the inclusion. Then $\pi_{1} \circ j: Y \rightarrow\{a\}$ and $\pi_{2} \circ j=\operatorname{Id}_{Y}$; in particular $j^{\star} E=\left(\pi_{1} \circ j\right)^{\star} \mathscr{O}(1)$ is the trivial bundle $Y \times \mathscr{O}(1)_{a}$. Consider now the composite morphism

$$
\begin{aligned}
H^{p, p}(Y, \mathbb{C}) \otimes H^{0}(P(V), \mathbb{O}(1)) & \longrightarrow H^{p, p}(X, E) \xrightarrow{j^{\star}} H^{p, p}(Y, \mathbb{C}) \otimes \mathbb{O}(1)_{a} \\
u \otimes s & \longmapsto \pi_{2}^{\star} u \otimes \pi_{1}^{\star} s,
\end{aligned}
$$

given by $u \otimes s \longmapsto\left(\pi_{2} \circ j\right)^{\star} u \otimes\left(\pi_{1} \circ j\right)^{\star} s=u \otimes s(a)$; it is surjective and $H^{p, p}(Y, \mathbb{C}) \neq 0$ for $0 \leqslant p \leqslant n-1$, so we have $H^{p, p}(X, E) \neq 0$.

## § 5. Vanishing Theorem for Partially Positive Line Bundles

Even in the case when the curvature form $\mathrm{i} \Theta(E)$ is not semi-positive, some cohomology groups of high tensor powers $E^{k}$ still vanish under suitable assumptions. The prototype of such results is the following assertion, which can be seen as a consequence of the Andreotti-Grauert theorem [Andreotti-Grauert 1962], see IX-?.?; the special case where $E$ is $>0$ (that is, $s=1$ ) is due to [Kodaira 1953] and [Serre 1956].
(5.1) Theorem. Let $F$ be a holomorphic vector bundle over a compact complex manifold $X$, s a positive integer and $E$ a hermitian line bundle such that $\mathrm{i} \Theta(E)$ has at least $n-s+1$ positive eigenvalues at every point $x \in X$. Then there exists an integer $k_{0} \geqslant 0$ such that

$$
H^{q}\left(X, E^{k} \otimes F\right)=0 \quad \text { for } \quad q \geqslant s \text { and } k \geqslant k_{0}
$$

Proof. The main idea is to construct a hermitian metric $\omega_{\varepsilon}$ on $X$ in such a way that all negative eigenvalues of $\mathrm{i} \Theta(E)$ with respect to $\omega_{\varepsilon}$ will be of small absolute value. Let $\omega$ denote a fixed hermitian metric on $X$ and let $\gamma_{1} \leqslant \ldots \leqslant \gamma_{n}$ be the corresponding eigenvalues of $\mathrm{i} \Theta(E)$.
(5.2) Lemma. Let $\psi \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$. If $A$ is a hermitian $n \times n$ matrix with eigenvalues $\lambda_{1} \leqslant \ldots \leqslant \lambda_{n}$ and corresponding eigenvectors $v_{1}, \ldots, v_{n}$, we define $\psi[A]$ as the hermitian matrix with eigenvalues $\psi\left(\lambda_{j}\right)$ and eigenvectors $v_{j}, 1 \leqslant j \leqslant n$. Then the map $A \longmapsto \psi[A]$ is $\mathscr{C}^{\infty}$ on $\operatorname{Herm}\left(\mathbb{C}^{n}\right)$.

Proof. Although the result is very well known, we give here a short proof. Without loss of generality, we may assume that $\psi$ is compactly supported. Then we have

$$
\psi[A]=\frac{1}{2 \pi} \int_{-\infty}^{+\infty} \widehat{\psi}(t) e^{i t A} d t
$$

where $\widehat{\psi}$ is the rapidly decreasing Fourier transform of $\psi$. The equality $\int_{0}^{t}(t-u)^{p} u^{q} d u=$ $p!q!/(p+q+1)$ ! and obvious power series developments yield

$$
D_{A}\left(e^{i t A}\right) \cdot B=i \int_{0}^{t} e^{\mathrm{i}(t-u) A} B e^{i u A} d u
$$

Since $e^{i u A}$ is unitary, we get $\left\|D_{A}\left(e^{i t A}\right)\right\| \leqslant|t|$. A differentiation under the integral sign and Leibniz' formula imply by induction on $k$ the bound $\left\|D_{A}^{k}\left(e^{i t A}\right)\right\| \leqslant|t|^{k}$. Hence $A \longmapsto \psi[A]$ is smooth.

Let us consider now the positive numbers

$$
t_{0}=\inf _{X} \gamma_{s}>0, \quad M=\sup _{X} \max _{j}\left|\gamma_{j}\right|>0
$$

We select a function $\psi_{\varepsilon} \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$ such that

$$
\psi_{\varepsilon}(t)=t \text { for } t \geqslant t_{0}, \quad \psi_{\varepsilon}(t) \geqslant t \text { for } 0 \leqslant t \leqslant t_{0}, \quad \psi_{\varepsilon}(t)=M / \varepsilon \text { for } t \leqslant 0
$$

By Lemma 5.2, $\omega_{\varepsilon}:=\psi_{\varepsilon}[\mathrm{i} \Theta(E)]$ is a smooth hermitian metric on $X$. Let us write

$$
\mathrm{i} \Theta(E)=\mathrm{i} \sum_{1 \leqslant j \leqslant n} \gamma_{j} \zeta_{j} \wedge \bar{\zeta}_{j}, \quad \omega_{\varepsilon}=\mathrm{i} \sum_{1 \leqslant j \leqslant n} \psi_{\varepsilon}\left(\gamma_{j}\right) \zeta_{j} \wedge \bar{\zeta}_{j}
$$

in an orthonormal basis $\left(\zeta_{1}, \ldots, \zeta_{n}\right)$ of $T^{\star} X$ for $\omega$. The eigenvalues of $\mathrm{i} \Theta(E)$ with respect to $\omega_{\varepsilon}$ are given by $\gamma_{j, \varepsilon}=\gamma_{j} / \psi_{\varepsilon}\left(\gamma_{j}\right)$ and the construction of $\psi_{\varepsilon}$ shows that $-\varepsilon \leqslant \gamma_{j, \varepsilon} \leqslant 1$, $1 \leqslant j \leqslant n$, and $\gamma_{j, \varepsilon}=1$ for $s \leqslant j \leqslant n$. Now, we have

$$
H^{q}\left(X, E^{k} \otimes F\right) \simeq H^{n, q}\left(X, E^{k} \otimes G\right)
$$

where $G=F \otimes K_{X}^{\star}$. Let $e,\left(g_{\lambda}\right)_{1 \leqslant \lambda \leqslant r}$ and $\left(\zeta_{j}\right)_{1 \leqslant j \leqslant n}$ denote orthonormal frames of $E$, $G$ and $\left(T^{\star} X, \omega_{\varepsilon}\right)$ respectively. For

$$
u=\sum_{|J|=q, \lambda} u_{J, \lambda} \zeta_{1} \wedge \ldots \wedge \zeta_{n} \wedge \bar{\zeta}_{J} \otimes e^{k} \otimes g_{\lambda} \in \Lambda^{n, q} T^{\star} X \otimes E^{k} \otimes G
$$

inequality (3.2) yields

$$
\left\langle\left[\mathrm{i} \Theta(E), \Lambda_{\varepsilon}\right] u, u\right\rangle_{\varepsilon}=\sum_{J, \lambda}\left(\sum_{j \in J} \gamma_{j, \varepsilon}\right)\left|u_{J, \lambda}\right|^{2} \geqslant(q-s+1-(s-1) \varepsilon)|u|^{2} .
$$

Choosing $\varepsilon=1 / s$ and $q \geqslant s$, the right hand side becomes $\geqslant(1 / s)|u|^{2}$. Since $\Theta\left(E^{k} \otimes G\right)=$ $k \Theta(E) \otimes \operatorname{Id}_{G}+\Theta(G)$, there exists an integer $k_{0}$ such that

$$
\left[\mathrm{i} \Theta\left(E^{k} \otimes G\right), \Lambda_{\varepsilon}\right]+T_{\omega_{\varepsilon}} \quad \text { acting on } \quad \Lambda^{n, q} T^{\star} X \otimes E^{k} \otimes G
$$

is positive definite for $q \geqslant s$ and $k \geqslant k_{0}$. The proof is complete.

## §6. Positivity Concepts for Vector Bundles

Let $E$ be a hermitian holomorphic vector bundle of rank $r$ over $X$, where $\operatorname{dim}_{\mathbb{C}} X=n$. Denote by $\left(e_{1}, \ldots, e_{r}\right)$ an orthonormal frame of $E$ over a coordinate patch $\Omega \subset X$ with complex coordinates $\left(z_{1}, \ldots, z_{n}\right)$, and

$$
\begin{equation*}
\mathrm{i} \Theta(E)=\mathrm{i} \sum_{1 \leqslant j, k \leqslant n, 1 \leqslant \lambda, \mu \leqslant r} c_{j k \lambda \mu} d z_{j} \wedge d \bar{z}_{k} \otimes e_{\lambda}^{\star} \otimes e_{\mu}, \quad \bar{c}_{j k \lambda \mu}=c_{k j \mu \lambda} \tag{6.1}
\end{equation*}
$$

the Chern curvature tensor. To $\mathrm{i} \Theta(E)$ corresponds a natural hermitian form $\theta_{E}$ on $T X \otimes E$ defined by

$$
\theta_{E}=\sum_{j, k, \lambda, \mu} c_{j k \lambda \mu}\left(d z_{j} \otimes e_{\lambda}^{\star}\right) \otimes\left(\overline{d z_{k} \otimes e_{\mu}^{\star}}\right)
$$

and such that

$$
\theta_{E}(u, u)=\sum_{j, k, \lambda, \mu} c_{j k \lambda \mu}(x) u_{j \lambda} \bar{u}_{k \mu}, \quad u \in T_{x} X \otimes E_{x} \cdot(6.2)
$$

(6.3) Definition ([Nakano 1955]). E is said to be Nakano positive (resp. Nakano seminegative) if $\theta_{E}$ is positive definite (resp. semi-negative) as a hermitian form on $T X \otimes E$, i.e. if for every $u \in T X \otimes E, u \neq 0$, we have

$$
\theta_{E}(u, u)>0 \quad(\text { resp } . \leqslant 0)
$$

We write $>_{\mathrm{Nak}}\left(\right.$ resp. $\left.\leqslant_{\mathrm{Nak}}\right)$ for Nakano positivity (resp. semi-negativity).
(6.4) Definition ([Griffiths 1969]). $E$ is said to be Griffiths positive (resp. Griffiths semi-negative) if for all $\xi \in T_{x} X, \xi \neq 0$ and $s \in E_{x}, s \neq 0$ we have

$$
\theta_{E}(\xi \otimes s, \xi \otimes s)>0 \quad(\text { resp } . \leqslant 0)
$$

We write $>_{\text {Grif }}\left(\right.$ resp.$\left.\leqslant_{\text {Grif }}\right)$ for Griffiths positivity (resp. semi-negativity).

It is clear that Nakano positivity implies Griffiths positivity and that both concepts coincide if $r=1$ (in the case of a line bundle, $E$ is merely said to be positive). One
can generalize further by introducing additional concepts of positivity which interpolate between Griffiths positivity and Nakano positivity.
(6.5) Definition. Let $T$ and $E$ be complex vector spaces of dimensions $n, r$ respectively, and let $\Theta$ be a hermitian form on $T \otimes E$.
a) A tensor $u \in T \otimes E$ is said to be of rank $m$ if $m$ is the smallest $\geqslant 0$ integer such that u can be written

$$
u=\sum_{j=1}^{m} \xi_{j} \otimes s_{j}, \quad \xi_{j} \in T, s_{j} \in E
$$

b) $\Theta$ is said to be m-positive (resp. m-semi-negative) if $\Theta(u, u)>0$ (resp. $\Theta(u, u) \leqslant 0)$ for every tensor $u \in T \otimes E$ of rank $\leqslant m, u \neq 0$. In this case, we write

$$
\Theta>_{m} 0 \quad\left(\text { resp. } \Theta \leqslant_{m} 0\right)
$$

We say that the bundle $E$ is $m$-positive if $\theta_{E}>_{m} 0$. Griffiths positivity corresponds to $m=1$ and Nakano positivity to $m \geqslant \min (n, r)$.
(6.6) Proposition. A bundle $E$ is Griffiths positive if and only if $E^{\star}$ is Griffiths negative.

Proof. By (V-4.3') we get $\mathrm{i} \Theta\left(E^{\star}\right)=-\mathrm{i} \Theta(E)^{\dagger}$, hence

$$
\theta_{E^{\star}}\left(\xi_{1} \otimes \bar{s}_{2}, \xi_{2} \otimes \bar{s}_{1}\right)=-\theta_{E}\left(\xi_{1} \otimes s_{1}, \xi_{2} \otimes s_{2}\right), \quad \forall \xi_{1}, \xi_{2} \in T X, \forall s_{1}, s_{2} \in E
$$

where $\bar{s}_{j}=\left\langle\bullet, s_{j}\right\rangle \in E^{\star}$. Proposition 6.6 follows immediately.
It should be observed that the corresponding duality property for Nakano positive bundles is not true. In fact, using (6.1) we get

$$
\mathrm{i} \Theta\left(E^{\star}\right)=-\mathrm{i} \sum_{j, k, \lambda, \mu} c_{j k \mu \lambda} d z_{j} \wedge d \bar{z}_{k} \otimes e_{\lambda}^{\star \star} \otimes e_{\mu}^{\star}
$$

$$
\begin{equation*}
\theta_{E^{\star}}(v, v)=-\sum_{j, k, \mu, \lambda} c_{j k \mu \lambda} v_{j \lambda} \bar{v}_{k \mu} \tag{6.7}
\end{equation*}
$$

for any $v=\sum v_{j \lambda}\left(\partial / \partial z_{j}\right) \otimes e_{\lambda}^{\star} \in T X \otimes E^{\star}$. The following example shows that Nakano positivity or negativity of $\theta_{E}$ and $\theta_{E^{\star}}$ are unrelated.
(6.8) Example. Let $H$ be the rank $n$ bundle over $\mathbb{P}^{n}$ defined in $\S \mathrm{V}-15$. For any $u=\sum u_{j \lambda}\left(\partial / \partial z_{j}\right) \otimes \widetilde{e}_{\lambda} \in T X \otimes H, v=\sum v_{j \lambda}\left(\partial / \partial z_{j}\right) \otimes \widetilde{e}_{\lambda}^{\star} \in T X \otimes H^{\star}, 1 \leqslant j, \lambda \leqslant n$, formula (V-15.9) implies

$$
\left\{\begin{align*}
\theta_{H}(u, u) & =\sum u_{j \lambda} \bar{u}_{\lambda j}  \tag{6.9}\\
\theta_{H^{\star}}(v, v) & =\sum v_{j j} \bar{v}_{\lambda \lambda}=\left|\sum v_{j j}\right|^{2}
\end{align*}\right.
$$

It is then clear that $H \geqslant_{\text {Grif }} 0$ and $H^{\star} \leqslant_{\text {Nak }} 0$, but $H$ is neither $\geqslant_{\text {Nak }} 0$ nor $\leqslant_{\text {Nak }} 0$.
(6.10) Proposition. Let $0 \rightarrow S \rightarrow E \rightarrow Q \rightarrow 0$ be an exact sequence of hermitian vector bundles. Then
a) $E \geqslant_{\text {Grif }} 0 \Longrightarrow Q \geqslant_{\text {Grif }} 0$,
b) $E \leqslant_{\text {Grif }} 0 \quad \Longrightarrow \quad S \leqslant_{\text {Grif }} 0$,
c) $E \leqslant_{\text {Nak }} 0 \quad \Longrightarrow \quad S \leqslant_{\text {Nak }} 0$,
and analogous implications hold true for strict positivity.
Proof. If $\beta$ is written $\sum d z_{j} \otimes \beta_{j}, \beta_{j} \in \operatorname{hom}(S, Q)$, then formulas (V-14.6) and (V-14.7) yield

$$
\begin{aligned}
& \mathrm{i} \Theta(S)=\mathrm{i} \Theta(E)_{\upharpoonright S}-\sum d z_{j} \wedge d \bar{z}_{k} \otimes \beta_{k}^{\star} \beta_{j}, \\
& \mathrm{i} \Theta(Q)=\mathrm{i} \Theta(E)_{\upharpoonright Q}+\sum d z_{j} \wedge d \bar{z}_{k} \otimes \beta_{j} \beta_{k}^{\star} .
\end{aligned}
$$

Since $\beta \cdot(\xi \otimes s)=\sum \xi_{j} \beta_{j} \cdot s$ and $\beta^{\star} \cdot(\xi \otimes s)=\sum \bar{\xi}_{k} \beta_{k}^{\star} \cdot s$ we get

$$
\begin{gathered}
\theta_{S}\left(\xi \otimes s, \xi^{\prime} \otimes s^{\prime}\right)=\theta_{E}\left(\xi \otimes s, \xi^{\prime} \otimes s^{\prime}\right)-\sum_{j, k} \xi_{j} \bar{\xi}_{k}^{\prime}\left\langle\beta_{j} \cdot s, \beta_{k} \cdot s^{\prime}\right\rangle, \\
\theta_{S}(u, u)=\theta_{E}(u, u)-|\beta \cdot u|^{2} \\
\theta_{Q}\left(\xi \otimes s, \xi^{\prime} \otimes s^{\prime}\right)=\theta_{E}\left(\xi \otimes s, \xi^{\prime} \otimes s^{\prime}\right)+\sum_{j, k} \xi_{j} \bar{\xi}_{k}^{\prime}\left\langle\beta_{k}^{\star} \cdot s, \beta_{j}^{\star} \cdot s^{\prime}\right\rangle \\
\theta_{Q}(\xi \otimes s, \xi \otimes s)=\theta_{E}(\xi \otimes s, \xi \otimes s)+\left|\beta^{\star} \cdot(\xi \otimes s)\right|^{2}
\end{gathered}
$$

Since $H$ is a quotient bundle of the trivial bundle $\underline{V}$, Example 6.8 shows that $E \geqslant_{\text {Nak }} 0$ does not imply $Q \geqslant_{\text {Nak }} 0$.

## § 7. Nakano Vanishing Theorem

Let $(X, \omega)$ be a compact Kähler manifold, $\operatorname{dim}_{\mathbb{C}} X=n$, and $E \longrightarrow X$ a hermitian vector bundle of rank $r$. We are going to compute explicitly the hermitian operator $[\mathrm{i} \Theta(E), \Lambda]$ acting on $\Lambda^{p, q} T^{\star} X \otimes E$. Let $x_{0} \in X$ and $\left(z_{1}, \ldots, z_{n}\right)$ be local coordinates such that $\left(\partial / \partial z_{1}, \ldots, \partial / \partial z_{n}\right)$ is an orthonormal basis of $(T X, \omega)$ at $x_{0}$. One can write

$$
\begin{aligned}
\omega_{x_{0}} & =\mathrm{i} \sum_{1 \leqslant j \leqslant n} d z_{j} \wedge d \bar{z}_{j}, \\
\mathrm{i} \Theta(E)_{x_{0}} & =\mathrm{i} \sum_{j, k, \lambda, \mu} c_{j k \lambda \mu} d z_{j} \wedge d \bar{z}_{k} \otimes e_{\lambda}^{\star} \otimes e_{\mu}
\end{aligned}
$$

where $\left(e_{1}, \ldots, e_{r}\right)$ is an orthonormal basis of $E_{x_{0}}$. Let

$$
u=\sum_{|J|=p,|K|=q, \lambda} u_{J, K, \lambda} d z_{J} \wedge d \bar{z}_{K} \otimes e_{\lambda} \in\left(\Lambda^{p, q} T^{\star} X \otimes E\right)_{x_{0}}
$$

A simple computation as in the proof of Prop. VI-8.3 gives

$$
\begin{aligned}
\Lambda u & \left.\left.=\mathrm{i}(-1)^{p} \sum_{J, K, \lambda, s} u_{J, K, \lambda}\left(\frac{\partial}{\partial z_{s}}\right\lrcorner d z_{J}\right) \wedge\left(\frac{\partial}{\partial \bar{z}_{s}}\right\lrcorner d \bar{z}_{K}\right) \otimes e_{\lambda}, \\
\mathrm{i} \Theta(E) \wedge u & =\mathrm{i}(-1)^{p} \sum_{j, k, \lambda, \mu, J, K} c_{j k \lambda \mu} u_{J, K, \lambda} d z_{j} \wedge d z_{J} \wedge d \bar{z}_{k} \wedge d \bar{z}_{K} \otimes e_{\mu} \\
{[\mathrm{i} \Theta(E), \Lambda] u } & \left.=\sum_{j, k, \lambda, \mu, J, K} c_{j k \lambda \mu} u_{J, K, \lambda} d z_{j} \wedge\left(\frac{\partial}{\partial z_{k}}\right\lrcorner d z_{J}\right) \wedge d \bar{z}_{K} \otimes e_{\mu} \\
& \left.+\sum_{j, k, \lambda, \mu, J, K} c_{j k \lambda \mu} u_{J, K, \lambda} d z_{J} \wedge d \bar{z}_{k} \wedge\left(\frac{\partial}{\partial \bar{z}_{j}}\right\lrcorner d \bar{z}_{K}\right) \otimes e_{\mu} \\
& -\sum_{j, \lambda, \mu, J, K} c_{j j \lambda \mu} u_{J, K, \lambda} d z_{J} \wedge d \bar{z}_{K} \otimes e_{\mu} .
\end{aligned}
$$

We extend the definition of $u_{J, K, \lambda}$ to non increasing multi-indices $J=\left(j_{s}\right), K=\left(k_{s}\right)$ by deciding that $u_{J, K, \lambda}=0$ if $J$ or $K$ contains identical components repeated and that $u_{J, K, \lambda}$ is alternate in the indices $\left(j_{s}\right),\left(k_{s}\right)$. Then the above equality can be written

$$
\begin{aligned}
\langle[\mathrm{i} \Theta(E), \Lambda] u, u\rangle & =\sum c_{j k \lambda \mu} u_{J, j S, \lambda} \bar{u}_{J, k S, \mu} \\
& +\sum c_{j k \lambda \mu} u_{k R, K, \lambda} \bar{u}_{j R, K, \mu} \\
& -\sum c_{j j \lambda \mu} u_{J, K, \lambda} \bar{u}_{J, K, \mu},
\end{aligned}
$$

extended over all indices $j, k, \lambda, \mu, J, K, R, S$ with $|R|=p-1,|S|=q-1$. This hermitian form appears rather difficult to handle for general $(p, q)$ because of sign compensation. Two interesting cases are $p=n$ and $q=n$.

- For $u=\sum u_{K, \lambda} d z_{1} \wedge \ldots \wedge d z_{n} \wedge d \bar{z}_{K} \otimes e_{\lambda}$ of type $(n, q)$, we get

$$
\begin{equation*}
\langle[\mathrm{i} \Theta(E), \Lambda] u, u\rangle=\sum_{|S|=q-1} \sum_{j, k, \lambda, \mu} c_{j k \lambda \mu} u_{j S, \lambda} \bar{u}_{k S, \mu}, \tag{7.1}
\end{equation*}
$$

because of the equality of the second and third summations in the general formula. Since $u_{j S, \lambda}=0$ for $j \in S$, the rank of the tensor $\left(u_{j S, \lambda}\right)_{j, \lambda} \in \mathbb{C}^{n} \otimes \mathbb{C}^{r}$ is in fact $\leqslant$ $\min \{n-q+1, r\}$. We obtain therefore:
(7.2) Lemma. Assume that $E>_{m} 0$ in the sense of Def. 6.5. Then the hermitian operator $[\mathrm{i} \Theta(E), \Lambda]$ is positive definite on $\Lambda^{n, q} T^{\star} X \otimes E$ for $q \geqslant 1$ and $m \geqslant \min \{n-q+$ $1, r\}$.
(7.3) Theorem. Let $X$ be a compact connected Kähler manifold of dimension $n$ and $E$ a hermitian vector bundle of rank $r$. If $\theta_{E} \geqslant_{m} 0$ on $X$ and $\theta_{E}>_{m} 0$ in at least one point, then

$$
H^{n, q}(X, E)=H^{q}\left(X, K_{X} \otimes E\right)=0 \quad \text { for } \quad q \geqslant 1 \quad \text { and } \quad m \geqslant \min \{n-q+1, r\} .
$$

- Similarly, for $u=\sum u_{J, \lambda} d z_{J} \wedge d \bar{z}_{1} \wedge \ldots \wedge d \bar{z}_{n} \otimes e_{\lambda}$ of type $(p, n)$, we get

$$
\langle[\mathrm{i} \Theta(E), \Lambda] u, u\rangle=\sum_{|R|=p-1} \sum_{j, k, \lambda, \mu} c_{j k \lambda \mu} u_{k R, \lambda} \bar{u}_{j R, \mu},
$$

because of the equality of the first and third summations in the general formula. The indices $j, k$ are twisted, thus $[\mathrm{i} \Theta(E), \Lambda]$ defines a positive hermitian form under the assumption $\mathrm{i} \Theta(E)^{\dagger}>_{m} 0$, i.e. $\mathrm{i} \Theta\left(E^{\star}\right)<_{m} 0$, with $m \geqslant \min \{n-p+1, r\}$. Serre duality $\left(H^{p, 0}(X, E)\right)^{\star}=H^{n-p, n}\left(X, E^{\star}\right)$ gives:
(7.4) Theorem. Let $X$ and $E$ be as above. If $\theta_{E} \leqslant_{m} 0$ on $X$ and $\theta_{E}<_{m} 0$ in at least one point, then

$$
H^{p, 0}(X, E)=H^{0}\left(X, \Omega_{X}^{p} \otimes E\right)=0 \quad \text { for } \quad p<n \quad \text { and } m \geqslant \min \{p+1, r\} .
$$

The special case $m=r$ yields:
(7.5) Corollary. For $X$ and $E$ as above:
a) Nakano vanishing theorem (1955):
$E \geqslant_{\text {Nak }} 0, \quad$ strictly in one point $\quad \Longrightarrow \quad H^{n, q}(X, E)=0 \quad$ for $q \geqslant 1$.
b) $E \leqslant \begin{aligned} & \text { Nak }\end{aligned} 0, \quad$ strictly in one point $\quad \Longrightarrow \quad H^{p, 0}(X, E)=0 \quad$ for $p<n$.

## § 8. Relations Between Nakano and Griffiths Positivity

It is clear that Nakano positivity implies Griffiths positivity. The main result of § 8 is the following "converse" to this property [Demailly-Skoda 1979].
(8.1) Theorem. For any hermitian vector bundle $E$,

$$
E>_{\text {Grif }} 0 \Longrightarrow E \otimes \operatorname{det} E>_{\text {Nak }} 0
$$

To prove this result, we first use (V-4.2') and (V-4.6). If $\operatorname{End}(E \otimes \operatorname{det} E)$ is identified to $\operatorname{hom}(E, E)$, one can write

$$
\begin{gathered}
\Theta(E \otimes \operatorname{det} E)=\Theta(E)+\operatorname{Tr}_{E}(\Theta(E)) \otimes \operatorname{Id}_{E}, \\
\theta_{E \otimes \operatorname{det} E}=\theta_{E}+\operatorname{Tr}_{E} \theta_{E} \otimes h,
\end{gathered}
$$

where $h$ denotes the hermitian metric on $E$ and where $\operatorname{Tr}_{E} \theta_{E}$ is the hermitian form on $T X$ defined by

$$
\operatorname{Tr}_{E} \theta_{E}(\xi, \xi)=\sum_{1 \leqslant \lambda \leqslant r} \theta_{E}\left(\xi \otimes e_{\lambda}, \xi \otimes e_{\lambda}\right), \xi \in T X
$$

for any orthonormal frame $\left(e_{1}, \ldots, e_{r}\right)$ of $E$. Theorem 8.1 is now a consequence of the following simple property of hermitian forms on a tensor product of complex vector spaces.
(8.2) Proposition. Let $T, E$ be complex vector spaces of respective dimensions n, $r$, and $h$ a hermitian metric on $E$. Then for every hermitian form $\Theta$ on $T \otimes E$

$$
\Theta>_{\text {Grif }} 0 \Longrightarrow \Theta+\operatorname{Tr}_{E} \Theta \otimes h>_{\text {Nak }} 0
$$

We first need a lemma analogous to Fourier inversion formula for discrete Fourier transforms.
(8.3) Lemma. Let $q$ be an integer $\geqslant 3$, and $x_{\lambda}, y_{\mu}, 1 \leqslant \lambda, \mu \leqslant r$, be complex numbers. Let $\sigma$ describe the set $U_{q}^{r}$ of $r$-tuples of $q$-th roots of unity and put

$$
x_{\sigma}^{\prime}=\sum_{1 \leqslant \lambda \leqslant r} x_{\lambda} \bar{\sigma}_{\lambda}, \quad y_{\sigma}^{\prime}=\sum_{1 \leqslant \mu \leqslant r} y_{\mu} \bar{\sigma}_{\mu}, \quad \sigma \in U_{q}^{r}
$$

Then for every pair $(\alpha, \beta), 1 \leqslant \alpha, \beta \leqslant r$, the following identity holds:

$$
q^{-r} \sum_{\sigma \in U_{q}^{r}} x_{\sigma}^{\prime} \bar{y}_{\sigma}^{\prime} \sigma_{\alpha} \bar{\sigma}_{\beta}= \begin{cases}x_{\alpha} \bar{y}_{\beta} & \text { if } \alpha \neq \beta, \\ \sum_{1 \leqslant \mu \leqslant r} x_{\mu} \bar{y}_{\mu} & \text { if } \alpha=\beta .\end{cases}
$$

Proof. The coefficient of $x_{\lambda} \bar{y}_{\mu}$ in the summation $q^{-r} \sum_{\sigma \in U_{q}^{r}} x_{\sigma}^{\prime} \bar{y}_{\sigma}^{\prime} \sigma_{\alpha} \bar{\sigma}_{\beta}$ is given by

$$
q^{-r} \sum_{\sigma \in U_{q}^{r}} \sigma_{\alpha} \bar{\sigma}_{\beta} \bar{\sigma}_{\lambda} \sigma_{\mu} .
$$

This coefficient equals 1 when the pairs $\{\alpha, \mu\}$ and $\{\beta, \lambda\}$ are equal (in which case $\sigma_{\alpha} \bar{\sigma}_{\beta} \bar{\sigma}_{\lambda} \sigma_{\mu}=1$ for any one of the $q^{r}$ elements of $\left.U_{q}^{r}\right)$. Hence, it is sufficient to prove that

$$
\sum_{\sigma \in U_{q}^{r}} \sigma_{\alpha} \bar{\sigma}_{\beta} \bar{\sigma}_{\lambda} \sigma_{\mu}=0
$$

when the pairs $\{\alpha, \mu\}$ and $\{\beta, \lambda\}$ are distinct.
If $\{\alpha, \mu\} \neq\{\beta, \lambda\}$, then one of the elements of one of the pairs does not belong to the other pair. As the four indices $\alpha, \beta, \lambda, \mu$ play the same role, we may suppose for example that $\alpha \notin\{\beta, \lambda\}$. Let us apply to $\sigma$ the substitution $\sigma \mapsto \tau$, where $\tau$ is defined by

$$
\tau_{\alpha}=e^{2 \pi \mathrm{i} / q} \sigma_{\alpha}, \tau_{\nu}=\sigma_{\nu} \quad \text { for } \quad \nu \neq \alpha .
$$

We get

$$
\sum_{\sigma} \sigma_{\alpha} \bar{\sigma}_{\beta} \bar{\sigma}_{\lambda} \sigma_{\mu}=\sum_{\tau}= \begin{cases}e^{2 \pi \mathrm{i} / q} \sum_{\sigma} & \text { if } \alpha \neq \mu \\ e^{4 \pi \mathrm{i} / q} \sum_{\sigma} & \text { if } \alpha=\mu\end{cases}
$$

Since $q \geqslant 3$ by hypothesis, it follows that

$$
\sum_{\sigma} \sigma_{\alpha} \bar{\sigma}_{\beta} \bar{\sigma}_{\lambda} \sigma_{\mu}=0 .
$$

Proof of Proposition 8.2.. Let $\left(t_{j}\right)_{1 \leqslant j \leqslant n}$ be a basis of $T,\left(e_{\lambda}\right)_{1 \leqslant \lambda \leqslant r}$ an orthonormal basis of $E$ and $\xi=\sum_{j} \xi_{j} t_{j} \in T, u=\sum_{j, \lambda} u_{j \lambda} t_{j} \otimes e_{\lambda} \in T \otimes E$. The coefficients $c_{j k \lambda \mu}$ of $\Theta$
with respect to the basis $t_{j} \otimes e_{\lambda}$ satisfy the symmetry relation $\bar{c}_{j k \lambda \mu}=c_{k j \mu \lambda}$, and we have the formulas

$$
\begin{aligned}
\Theta(u, u) & =\sum_{j, k, \lambda, \mu} c_{j k \lambda \mu} u_{j \lambda} \bar{u}_{k \mu}, \\
\operatorname{Tr}_{E} \Theta(\xi, \xi) & =\sum_{j, k, \lambda} c_{j k \lambda \lambda} \xi_{j} \bar{\xi}_{k}, \\
\left(\Theta+\operatorname{Tr}_{E} \Theta \otimes h\right)(u, u) & =\sum_{j, k, \lambda, \mu} c_{j k \lambda \mu} u_{j \lambda} \bar{u}_{k \mu}+c_{j k \lambda \lambda} u_{j \mu} \bar{u}_{k \mu} .
\end{aligned}
$$

For every $\sigma \in U_{q}^{r}$ (cf. Lemma 8.3), put

$$
\begin{aligned}
u_{j \sigma}^{\prime} & =\sum_{1 \leqslant \lambda \leqslant r} u_{j \lambda} \bar{\sigma}_{\lambda} \in \mathbb{C} \\
\widehat{u}_{\sigma} & =\sum_{j} u_{j \sigma}^{\prime} t_{j} \in T \quad, \quad \widehat{e}_{\sigma}=\sum_{\lambda} \sigma_{\lambda} e_{\lambda} \in E .
\end{aligned}
$$

Lemma 8.3 implies

$$
\begin{aligned}
q^{-r} \sum_{\sigma \in U_{q}^{r}} \Theta\left(\widehat{u}_{\sigma} \otimes \widehat{e}_{\sigma}, \widehat{u}_{\sigma} \otimes \widehat{e}_{\sigma}\right) & =q^{-r} \sum_{\sigma \in U_{q}^{r}} c_{j k \lambda \mu} u_{j \sigma}^{\prime} \bar{u}_{k \sigma}^{\prime} \sigma_{\lambda} \bar{\sigma}_{\mu} \\
& =\sum_{j, k, \lambda \neq \mu} c_{j k \lambda \mu} u_{j \lambda} \bar{u}_{k \mu}+\sum_{j, k, \lambda, \mu} c_{j k \lambda \lambda} u_{j \mu} \bar{u}_{k \mu}
\end{aligned}
$$

The Griffiths positivity assumption shows that the left hand side is $\geqslant 0$, hence

$$
\left(\Theta+\operatorname{Tr}_{E} \Theta \otimes h\right)(u, u) \geqslant \sum_{j, k, \lambda} c_{j k \lambda \lambda} u_{j \lambda} \bar{u}_{k \lambda} \geqslant 0
$$

with strict positivity if $\Theta>_{\text {Grif }} 0$ and $u \neq 0$.
(8.4) Example. Take $E=H$ over $\mathbb{P}^{n}=P(V)$. The exact sequence

$$
0 \longrightarrow \mathbb{O}(-1) \longrightarrow \underline{V} \longrightarrow H \longrightarrow 0
$$

implies $\operatorname{det} \underline{V}=\operatorname{det} H \otimes \mathscr{O}(-1)$. Since $\operatorname{det} \underline{V}$ is a trivial bundle, we get (non canonical) isomorphisms

$$
\begin{aligned}
\operatorname{det} H & \simeq \mathbb{O}(1), \\
T \mathbb{P}^{n} & =H \otimes \overparen{O}(1) \simeq H \otimes \operatorname{det} H .
\end{aligned}
$$

We already know that $H \geqslant_{\text {Grif }} 0$, hence $T \mathbb{P}^{n} \geqslant_{\text {Nak }} 0$. A direct computation based on (6.9) shows that

$$
\begin{aligned}
\theta_{T \mathbb{P}^{n}}(u, u) & =\left(\theta_{H}+\operatorname{Tr}_{H} \theta_{H} \otimes h\right)(u, u) \\
& =\sum_{1 \leqslant j, k \leqslant n} u_{j k} \bar{u}_{k j}+u_{j k} \bar{u}_{j k}=\frac{1}{2} \sum_{1 \leqslant j, k \leqslant n}\left|u_{j k}+u_{k j}\right|^{2} .
\end{aligned}
$$

In addition, we have $T \mathbb{P}^{n}>_{\text {Grif }} 0$. However, the Serre duality theorem gives

$$
\begin{aligned}
H^{q}\left(\mathbb{P}^{n}, K_{\mathbb{P}^{n}} \otimes T \mathbb{P}^{n}\right)^{\star} & \simeq H^{n-q}\left(\mathbb{P}^{n}, T^{\star} \mathbb{P}^{n}\right) \\
& =H^{1, n-q}\left(\mathbb{P}^{n}, \mathbb{C}\right)= \begin{cases}\mathbb{C} & \text { if } q=n-1, \\
0 & \text { if } q \neq n-1\end{cases}
\end{aligned}
$$

For $n \geqslant 2$, Th. 7.3 implies that $T \mathbb{P}^{n}$ has no hermitian metric such that $\theta_{T \mathbb{P}^{n}} \geqslant_{2} 0$ on $\mathbb{P}^{n}$ and $\theta_{T \mathbb{P}^{n}}>_{2} 0$ in one point. This shows that the notion of 2-positivity is actually stronger than 1-positivity (i.e. Griffiths positivity).
(8.5) Remark. Since $\operatorname{Tr}_{H} \theta_{H}=\theta_{\sigma(1)}$ is positive and $\theta_{T \mathbb{P}^{n}}$ is not $>_{\text {Nak }} 0$ when $n \geqslant 2$, we see that Prop. 8.2 is best possible in the sense that there cannot exist any constant $c<1$ such that

$$
\Theta>_{\text {Grif }} 0 \Longrightarrow \Theta+c \operatorname{Tr}_{E} \Theta \otimes h \geqslant_{\text {Nak }} 0
$$

## § 9. Applications to Griffiths Positive Bundles

We first need a preliminary result.
(9.1) Proposition. Let $T$ be a complex vector space and ( $E, h$ ) a hermitian vector space of respective dimensions $n, r$ with $r \geqslant 2$. Then for any hermitian form $\Theta$ on $T \otimes E$ and any integer $m \geqslant 1$

$$
\Theta>_{\text {Grif }} 0 \quad \Longrightarrow \quad m \operatorname{Tr}_{E} \Theta \otimes h-\Theta>_{m} 0
$$

Proof. Let us distinguish two cases.
a) $m=1$. Let $u \in T \otimes E$ be a tensor of rank 1 . Then $u$ can be written $u=\xi_{1} \otimes e_{1}$ with $\xi_{1} \in T, \xi_{1} \neq 0$, and $e_{1} \in E,\left|e_{1}\right|=1$. Complete $e_{1}$ into an orthonormal basis $\left(e_{1}, \ldots, e_{r}\right)$ of $E$. One gets immediately

$$
\begin{aligned}
\left(\operatorname{Tr}_{E} \Theta \otimes h\right)(u, u) & =\operatorname{Tr}_{E} \Theta\left(\xi_{1}, \xi_{1}\right)=\sum_{1 \leqslant \lambda \leqslant r} \Theta\left(\xi_{1} \otimes e_{\lambda}, \xi_{1} \otimes e_{\lambda}\right) \\
& >\Theta\left(\xi_{1} \otimes e_{1}, \xi_{1} \otimes e_{1}\right)=\Theta(u, u)
\end{aligned}
$$

b) $m \geqslant 2$. Every tensor $u \in T \otimes E$ of rank $\leqslant m$ can be written

$$
u=\sum_{1 \leqslant \lambda \leqslant q} \xi_{\lambda} \otimes e_{\lambda} \quad, \quad \xi_{\lambda} \in T
$$

with $q=\min (m, r)$ and $\left(e_{\lambda}\right)_{1 \leqslant \lambda \leqslant r}$ an orthonormal basis of $E$. Let $F$ be the vector subspace of $E$ generated by $\left(e_{1}, \ldots, e_{q}\right)$ and $\Theta_{F}$ the restriction of $\Theta$ to $T \otimes F$. The first part shows that

$$
\Theta^{\prime}:=\operatorname{Tr}_{F} \Theta_{F} \otimes h-\Theta_{F}>_{\text {Grif }} 0 .
$$

Proposition 9.2 applied to $\Theta^{\prime}$ on $T \otimes F$ yields

$$
\Theta^{\prime}+\operatorname{Tr}_{F} \Theta^{\prime} \otimes h=q \operatorname{Tr}_{F} \Theta_{F} \otimes h-\Theta_{F}>_{q} 0 .
$$

Since $u \in T \otimes F$ is of rank $\leqslant q \leqslant m$, we get (for $u \neq 0$ )

$$
\begin{aligned}
\Theta(u, u)=\Theta_{F}(u, u) & <q\left(\operatorname{Tr}_{F} \Theta_{F} \otimes h\right)(u, u) \\
& =q \sum_{1 \leqslant j, \lambda \leqslant q} \Theta\left(\xi_{j} \otimes e_{\lambda}, \xi_{j} \otimes e_{\lambda}\right) \leqslant m \operatorname{Tr}_{E} \Theta \otimes h(u, u) .
\end{aligned}
$$

Proposition 9.1 is of course also true in the semi-positive case. From these facts, we deduce
(9.2) Theorem. Let $E$ be a Griffiths (semi-)positive bundle of rank $r \geqslant 2$. Then for any integer $m \geqslant 1$

$$
E^{\star} \otimes(\operatorname{det} E)^{m}>_{m} 0 \quad\left(\text { resp } . \quad \geqslant_{m} 0\right)
$$

Proof. Apply Prop. 8.1 to $\Theta=-\theta_{E^{\star}}>_{\text {Grif }} 0$ and observe that

$$
\theta_{\operatorname{det} E}=-\theta_{\operatorname{det} E^{\star}}=\operatorname{Tr}_{E^{\star}} \Theta
$$

(9.3) Theorem. Let $0 \rightarrow S \rightarrow E \rightarrow Q \rightarrow 0$ be an exact sequence of hermitian vector bundles. Then for any $m \geqslant 1$

$$
E>_{m} 0 \quad \Longrightarrow \quad S \otimes(\operatorname{det} Q)^{m}>_{m} 0 .
$$

Proof. Formulas (V-14.6) and (V-14.7) imply

$$
\begin{aligned}
& \mathrm{i} \Theta(S)>_{m} \mathrm{i} \beta^{\star} \wedge \beta \quad, \quad \mathrm{i} \Theta(Q)>_{m} \mathrm{i} \beta \wedge \beta^{\star} \\
& \mathrm{i} \Theta(\operatorname{det} Q)=\operatorname{Tr}_{Q}(\mathrm{i} \Theta(Q))>\operatorname{Tr}_{Q}\left(\mathrm{i} \beta \wedge \beta^{\star}\right)
\end{aligned}
$$

If we write $\beta=\sum d z_{j} \otimes \beta_{j}$ as in the proof of Prop. 6.10, then

$$
\begin{aligned}
\operatorname{Tr}_{Q}\left(\mathrm{i} \beta \wedge \beta^{\star}\right) & =\sum i d z_{j} \wedge d \bar{z}_{k} \operatorname{Tr}_{Q}\left(\beta_{j} \beta_{k}^{\star}\right) \\
& =\sum i d z_{j} \wedge d \bar{z}_{k} \operatorname{Tr}_{S}\left(\beta_{k}^{\star} \beta_{j}\right)=\operatorname{Tr}_{S}\left(-\mathrm{i} \beta^{\star} \wedge \beta\right)
\end{aligned}
$$

Furthermore, it has been already proved that $-\mathrm{i} \beta^{\star} \wedge \beta \geqslant_{\text {Nak }} 0$. By Prop. 8.1 applied to the corresponding hermitian form $\Theta$ on $T X \otimes S$, we get

$$
m \operatorname{Tr}_{S}\left(-\mathrm{i} \beta^{\star} \wedge \beta\right) \otimes \operatorname{Id}_{S}+\mathrm{i} \beta^{\star} \wedge \beta \geqslant_{m} 0
$$

and Th. 9.3 follows.
(9.4) Corollary. Let $X$ be a compact $n$-dimensional complex manifold, $E$ a vector bundle of rank $r \geqslant 2$ and $m \geqslant 1$ an integer. Then
a) $E>_{\text {Grif }} 0 \Longrightarrow H^{n, q}(X, E \otimes \operatorname{det} E)=0 \quad$ for $q \geqslant 1$;
b) $E>_{\text {Grif }} 0 \Longrightarrow H^{n, q}\left(X, E^{\star} \otimes(\operatorname{det} E)^{m}\right)=0 \quad$ for $q \geqslant 1$ and $m \geqslant \min \{n-q+1, r\}$;
c) Let $0 \rightarrow S \rightarrow E \rightarrow Q \rightarrow 0$ be an exact sequence of vector bundles and $m=\min \{n-$ $q+1, \operatorname{rk} S\}, q \geqslant 1$. If $E>_{m} 0$ and if $L$ is a line bundle such that $L \otimes(\operatorname{det} Q)^{-m} \geqslant 0$, then

$$
H^{n, q}(X, S \otimes L)=0
$$

Proof. Immediate consequence of Theorems 7.3, 8.1, 9.2 and 9.3.
Note that under our hypotheses $\omega=\mathrm{i} \operatorname{Tr}_{E} \Theta(E)=\mathrm{i} \Theta\left(\Lambda^{r} E\right)$ is a Kähler metric on $X$. Corollary 2.5 shows that it is enough in a), b), c) to assume semi-positivity and strict positivity in one point (this is true a priori only if $X$ is supposed in addition to be Kähler, but this hypothesis can be removed by means of Siu's result mentioned after (4.5).
a) is in fact a special case of a result of [Griffiths 1969], which we will prove in full generality in volume II (see the chapter on vanishing theorems for ample vector bundles); property b) will be also considerably strengthened there. Property c) is due to [Skoda 1978] for $q=0$ and to [Demailly 1982c] in general. Let us take the tensor product of the exact sequence in c ) with $(\operatorname{det} Q)^{l}$. The corresponding long cohomology exact sequence implies that the natural morphism

$$
H^{n, q}\left(X, E \otimes(\operatorname{det} Q)^{l}\right) \longrightarrow H^{n, q}\left(X, Q \otimes(\operatorname{det} Q)^{l}\right)
$$

is surjective for $q \geqslant 0$ and $l, m \geqslant \min \{n-q$, rk $S\}$, bijective for $q \geqslant 1$ and $l, m \geqslant$ $\min \{n-q+1, \mathrm{rk} S\}$.

## $\S$ 10. Cohomology Groups of $\mathfrak{G}(k)$ over $\mathbb{P}^{n}$

As an illustration of the above results, we compute now the cohomology groups of all line bundles $\mathfrak{O}(k) \rightarrow \mathbb{P}^{n}$. This precise evaluation will be needed in the proof of a general vanishing theorem for vector bundles, due to Le Potier (see volume II). As in §V-15, we consider a complex vector space $V$ of dimension $n+1$ and the exact sequence

$$
\begin{equation*}
0 \longrightarrow \mathbb{O}(-1) \longrightarrow \underline{V} \longrightarrow H \longrightarrow 0 \tag{10.1}
\end{equation*}
$$

of vector bundles over $\mathbb{P}^{n}=P(V)$. We thus have $\operatorname{det} \underline{V}=\operatorname{det} H \otimes \mathbb{O}(-1)$, and as $T P(V)=H \otimes \mathscr{O}(1)$ by Th. V-15.7, we find

$$
\begin{equation*}
K_{P(V)}=\operatorname{det} T^{\star} P(V)=\operatorname{det} H^{\star} \otimes \mathscr{O}(-n)=\operatorname{det} \underline{V}^{\star} \otimes \mathscr{O}(-n-1) \tag{10.2}
\end{equation*}
$$

where $\operatorname{det} \underline{V}$ is a trivial line bundle.
Before going further, we need some notations. For every integer $k \in \mathbb{N}$, we consider the homological complex $C^{\bullet}, k\left(V^{\star}\right)$ with differential $\gamma$ such that

$$
\left\{\begin{array}{l}
C^{p, k}\left(V^{\star}\right)=\Lambda^{p} V^{\star} \otimes S^{k-p} V^{\star}, \quad 0 \leqslant p \leqslant k  \tag{10.3}\\
\quad=0 \quad \text { otherwise }, \\
\gamma: \quad \Lambda^{p} V^{\star} \otimes S^{k-p} V^{\star} \longrightarrow \Lambda^{p-1} V^{\star} \otimes S^{k-p+1} V^{\star}
\end{array}\right.
$$

where $\gamma$ is the linear map obtained by contraction with the Euler vector field $\operatorname{Id}_{V} \in V \otimes$ $V^{\star}$, through the obvious maps $V \otimes \Lambda^{p} V^{\star} \longrightarrow \Lambda^{p-1} V^{\star}$ and $V^{\star} \otimes S^{k-p} V^{\star} \longrightarrow S^{k-p+1} V^{\star}$. If $\left(z_{0}, \ldots, z_{n}\right)$ are coordinates on $V$, the module $C^{p, k}\left(V^{\star}\right)$ can be identified with the space of $p$-forms

$$
\alpha(z)=\sum_{|I|=p} \alpha_{I}(z) d z_{I}
$$

where the $\alpha_{I}$ 's are homogeneous polynomials of degree $k-p$. The differential $\gamma$ is given by contraction with the Euler vector field $\xi=\sum_{0 \leqslant j \leqslant n} z_{j} \partial / \partial z_{j}$.

Let us denote by $Z^{p, k}\left(V^{\star}\right)$ the space of $p$-cycles of $C^{\bullet, k}\left(V^{\star}\right)$, i.e. the space of forms $\alpha \in C^{p, k}\left(V^{\star}\right)$ such that $\left.\xi\right\lrcorner \alpha=0$. The exterior derivative $d$ also acts on $C^{\bullet, k}\left(V^{\star}\right)$; we have

$$
d: \quad C^{p, k}\left(V^{\star}\right) \longrightarrow C^{p+1, k}\left(V^{\star}\right)
$$

and a trivial computation shows that $d \gamma+\gamma d=k \cdot \operatorname{Id}_{C} \bullet{ }^{k}\left(V^{\star}\right)$.
(10.4) Theorem. For $k \neq 0, C^{\bullet, k}\left(V^{\star}\right)$ is exact and there exist canonical isomorphisms

$$
C^{\bullet, k}\left(V^{\star}\right)=\Lambda^{p} V^{\star} \otimes S^{k-p} V^{\star} \simeq Z^{p, k}\left(V^{\star}\right) \oplus Z^{p-1, k}\left(V^{\star}\right)
$$

Proof. The identity $d \gamma+\gamma d=k$. Id implies the exactness. The isomorphism is given by $\frac{1}{k} \gamma d \oplus \gamma$ and its inverse by $\mathscr{P}_{1}+\frac{1}{k} d \circ \mathscr{P}_{2}$.

Let us consider now the canonical mappings

$$
\pi: V \backslash\{0\} \longrightarrow P(V), \quad \mu^{\prime}: V \backslash\{0\} \longrightarrow \mathbb{O}(-1)
$$

defined in $\S V-15$. As $T_{[z]} P(V) \simeq V / \mathbb{C} \xi(z)$ for all $z \in V \backslash\{0\}$, every form $\alpha \in Z^{p, k}\left(V^{\star}\right)$ defines a holomorphic section of $\pi^{\star}\left(\Lambda^{p} T^{\star} P(V)\right), \alpha(z)$ being homogeneous of degree $k$ with respect to $z$. Hence $\alpha(z) \otimes \mu^{\prime}(z)^{-k}$ is a holomorphic section of $\pi^{\star}\left(\Lambda^{p} T^{\star} P(V) \otimes \mathscr{O}(k)\right)$, and since its homogeneity degree is 0 , it induces a holomorphic section of $\Lambda^{p} T^{\star} P(V) \otimes$ $\mathfrak{O}(k)$. We thus have an injective morphism

$$
\begin{equation*}
Z^{p, k}\left(V^{\star}\right) \longrightarrow H^{p, 0}(P(V), \mathscr{O}(k)) \tag{10.5}
\end{equation*}
$$

(10.6) Theorem. The groups $H^{p, 0}(P(V), \mathscr{O}(k))$ are given by
a) $H^{p, 0}(P(V), \mathscr{O}(k)) \simeq Z^{p, k}\left(V^{\star}\right)$ for $k \geqslant p \geqslant 0$,
b) $H^{p, 0}(P(V), \mathscr{O}(k))=0 \quad$ for $k \leqslant p$ and $(k, p) \neq(0,0)$.

Proof. Let $s$ be a holomorphic section of $\Lambda^{p} T^{\star} P(V) \otimes \mathscr{O}(k)$. Set

$$
\alpha(z)=\left(d \pi_{z}\right)^{\star}\left(s([z]) \otimes \mu^{\prime}(z)^{k}\right), \quad z \in V \backslash\{0\}
$$

Then $\alpha$ is a holomorphic $p$-form on $V \backslash\{0\}$ such that $\xi\lrcorner \alpha=0$, and the coefficients of $\alpha$ are homogeneous of degree $k-p$ on $V \backslash\{0\}$ (recall that $d \pi_{\lambda z}=\lambda^{-1} d \pi_{z}$ ). It follows that $\alpha=0$ if $k<p$ and that $\alpha \in Z^{p, k}\left(V^{\star}\right)$ if $k \geqslant p$. The injective morphism (10.5) is therefore also surjective. Finally, $Z^{p, p}\left(V^{\star}\right)=0$ for $p=k \neq 0$, because of the exactness of $C^{\bullet}, k\left(V^{\star}\right)$ when $k \neq 0$. The proof is complete.
(10.7) Theorem. The cohomology groups $H^{p, q}(P(V), \mathscr{O}(k))$ vanish in the following cases:
a) $q \neq 0, n, p$;
b) $q=0, \quad k \leqslant p$ and $(k, p) \neq(0,0)$;
c) $q=n, \quad k \geqslant-n+p$ and $(k, p) \neq(0, n)$;
d) $q=p \neq 0, n, \quad k \neq 0$.

The remaining non vanishing groups are:
Б) $H^{p, 0}(P(V), \overparen{\odot}(k)) \simeq Z^{p, k}\left(V^{\star}\right)$ for $k>p ;$
$\overline{\mathrm{c}}) H^{p, n}(P(V), \mathscr{O}(k)) \simeq Z^{n-p,-k}(V)$ for $k<-n+p$;
$\overline{\mathrm{d}}) H^{p, p}(P(V), \mathbb{C})=\mathbb{C}, \quad 0 \leqslant p \leqslant n$.
Proof. • $\overline{\mathrm{d}}$ ) is already known, and so is a) when $k=0$ (Th. VI-13.3).

- b) and $\bar{b}$ ) follow from Th. 10.6, and c), $\bar{c}$ ) are equivalent to b), $\bar{b}$ ) via Serre duality:

$$
H^{p, q}(P(V), \mathscr{O}(k))^{\star}=H^{n-p, n-q}(P(V), \mathbb{O}(-k))
$$

thanks to the canonical isomorphism $\left(Z^{p, k}(V)\right)^{\star}=Z^{p, k}\left(V^{\star}\right)$.

- Let us prove now property a) when $k \neq 0$ and property d). By Serre duality, we may assume $k>0$. Then

$$
\Lambda^{p} T^{\star} P(V) \simeq K_{P(V)} \otimes \Lambda^{n-p} T P(V)
$$

It is very easy to verify that $E \geqslant_{\text {Nak }} 0$ implies $\Lambda^{s} E \geqslant_{\text {Nak }} 0$ for every integer $s$. Since $T P(V) \geqslant_{\text {Nak }} 0$, we get therefore

$$
F=\Lambda^{n-p} T P(V) \otimes \mathscr{O}(k)>_{\text {Nak }} 0 \quad \text { for } \quad k>0,
$$

and the Nakano vanishing theorem implies

$$
\begin{aligned}
H^{p, q}(P(V), \mathscr{O}(k)) & =H^{q}\left(P(V), \Lambda^{p} T^{\star} P(V) \otimes \mathscr{O}(k)\right) \\
& =H^{q}\left(P(V), K_{P(V)} \otimes F\right)=0, \quad q \geqslant 1 .
\end{aligned}
$$

## § 11. Ample Vector Bundles

## §11.A. Globally Generated Vector Bundles

All definitions concerning ampleness are purely algebraic and do not involve differential geometry. We shall see however that ampleness is intimately connected with the differential geometric notion of positivity. For a general discussion of properties of ample vector bundles in arbitrary characteristic, we refer to [Hartshorne 1966].
(11.1) Definition. Let $E \rightarrow X$ be a holomorphic vector bundle over an arbitrary complex manifold $X$.
a) $E$ is said to be globally generated if for every $x \in X$ the evaluation map $H^{0}(X, E) \rightarrow$ $E_{x}$ is onto.
b) $E$ is said to be semi-ample if there exists an integer $k_{0}$ such that $S^{k} E$ is globally generated for $k \geqslant k_{0}$.

Any quotient of a trivial vector bundle is globally generated, for example the tautological quotient vector bundle $Q$ over the Grassmannian $G_{r}(V)$ is globally generated. Conversely, every globally generated vector bundle $E$ of rank $r$ is isomorphic to the quotient of a trivial vector bundle of rank $\leqslant n+r$, as shown by the following result.
(11.2) Proposition. If a vector bundle $E$ of rank $r$ is globally generated, then there exists a finite dimensional subspace $V \subset H^{0}(X, E)$, $\operatorname{dim} V \leqslant n+r$, such that $V$ generates all fibers $E_{x}, x \in X$.

Proof. Put an arbitrary hermitian metric on $E$ and consider the Fréchet space $\mathscr{F}=$ $H^{0}(X, E)^{\oplus n+r}$ of $(n+r)$-tuples of holomorphic sections of $E$, endowed with the topology of uniform convergence on compact subsets of $X$. For every compact set $K \subset X$, we set

$$
A(K)=\left\{\left(s_{1}, \ldots, s_{n+r}\right) \in \mathscr{F} \text { which do not generate } E \text { on } K\right\} .
$$

It is enough to prove that $A(K)$ is of first category in $\mathscr{F}$ : indeed, Baire's theorem will imply that $A(X)=\bigcup A\left(K_{\nu}\right)$ is also of first category, if $\left(K_{\nu}\right)$ is an exhaustive sequence of compact subsets of $X$. It is clear that $A(K)$ is closed, because $A(K)$ is characterized by the closed condition

$$
\min _{K} \sum_{i_{1}<\cdots<i_{r}}\left|s_{i_{1}} \wedge \cdots \wedge s_{i_{r}}\right|=0
$$

It is therefore sufficient to prove that $A(K)$ has no interior point. By hypothesis, each fiber $E_{x}, x \in K$, is generated by $r$ global sections $s_{1}^{\prime}, \ldots, s_{r}^{\prime}$. We have in fact $s_{1}^{\prime} \wedge \cdots \wedge s_{r}^{\prime} \neq$ 0 in a neighborhood $U_{x}$ of $x$. By compactness of $K$, there exist finitely many sections $s_{1}^{\prime}, \ldots, s_{N}^{\prime}$ which generate $E$ in a neighborhood $\Omega$ of the set $K$.

If $T$ is a complex vector space of dimension $r$, define $R_{k}\left(T^{p}\right)$ as the set of $p$-tuples $\left(x_{1}, \ldots, x_{p}\right) \in T^{p}$ of rank $k$. Given $a \in R_{k}\left(T^{p}\right)$, we can reorder the $p$-tuple in such a way that $a_{1} \wedge \cdots \wedge a_{k} \neq 0$. Complete these $k$ vectors into a basis $\left(a_{1}, \ldots, a_{k}, b_{1}, \ldots, b_{r-k}\right)$ of $T$. For every point $x \in T^{p}$ in a neighborhood of $a$, then $\left(x_{1}, \ldots, x_{k}, b_{1}, \ldots, b_{r-k}\right)$ is again a basis of $T$. Therefore, we will have $x \in R_{k}\left(T^{p}\right)$ if and only if the coordinates of $x_{l}, k+1 \leqslant l \leqslant N$, relative to $b_{1}, \ldots, b_{r-k}$ vanish. It follows that $R_{k}\left(T^{p}\right)$ is a (non closed) submanifold of $T^{p}$ of codimension $(r-k)(p-k)$.

Now, we have a surjective affine bundle-morphism

$$
\begin{aligned}
\Phi: \Omega \times \mathbb{C}^{N(n+r)} & \longrightarrow E^{n+r} \\
(x, \lambda) & \longmapsto\left(s_{j}(x)+\sum_{1 \leqslant k \leqslant N} \lambda_{j k} s_{k}^{\prime}(x)\right)_{1 \leqslant j \leqslant n+r} .
\end{aligned}
$$

Therefore $\Phi^{-1}\left(R_{k}\left(E^{n+r}\right)\right)$ is a locally trivial differentiable bundle over $\Omega$, and the codimension of its fibers in $\mathbb{C}^{N(n+r)}$ is $(r-k)(n+r-k) \geqslant n+1$ if $k<r$; it follows that the dimension of the total space $\Phi^{-1}\left(R_{k}\left(E^{n+r}\right)\right)$ is $\leqslant N(n+r)-1$. By Sard's theorem

$$
\bigcup_{k<r} \mathscr{P}_{2}\left(\Phi^{-1}\left(R_{k}\left(E^{n+r}\right)\right)\right)
$$

is of zero measure in $\mathbb{C}^{N(n+r)}$. This means that for almost every value of the parameter $\lambda$ the vectors $s_{j}(x)+\sum_{k} \lambda_{j k} s_{k}^{\prime}(x) \in E_{x}, 1 \leqslant j \leqslant n+r$, are of maximum rank $r$ at each point $x \in \Omega$. Therefore $A(K)$ has no interior point.

Assume now that $V \subset H^{0}(X, E)$ generates $E$ on $X$. Then there is an exact sequence

$$
\begin{equation*}
0 \longrightarrow S \longrightarrow \underline{V} \longrightarrow E \longrightarrow 0 \tag{11.3}
\end{equation*}
$$

of vector bundles over $X$, where $S_{x}=\{s \in V ; s(x)=0\}, \operatorname{codim}_{V} S_{x}=r$. One obtains therefore a commutative diagram

where $\psi_{V}, \Psi_{V}$ are the holomorphic maps defined by

$$
\begin{aligned}
& \psi_{V}(x)=S_{x}, \quad x \in X \\
& \Psi_{V}(u)=\{s \in V ; s(x)=u\} \in V / S_{x}, \quad u \in E_{x}
\end{aligned}
$$

In particular, we see that every globally generated vector bundle $E$ of rank $r$ is the pullback of the tautological quotient vector bundle $Q$ of rank $r$ over the Grassmannian by means of some holomorphic map $X \longrightarrow G_{r}(V)$. In the special case when $\operatorname{rk} E=r=1$, the above diagram becomes

$$
\begin{array}{ccc}
E & \xrightarrow{\Psi_{V}} & \mathbb{O}(1) \\
\downarrow & & \downarrow \\
X & \xrightarrow{\psi_{V}} & P\left(V^{\star}\right)
\end{array}
$$

(11.5) Corollary. If $E$ is globally generated, then $E$ possesses a hermitian metric such that $E \geqslant_{\text {Grif }} 0$ (and also $E^{\star} \leqslant_{\text {Nak }} 0$ ).

Proof. Apply Prop. 6.11 to the exact sequence (11.3), where $\underline{V}$ is endowed with an arbitrary hermitian metric.

When $E$ is of rank $r=1$, then $S^{k} E=E^{\otimes k}$ and any hermitian metric of $E^{\otimes k}$ yields a metric on $E$ after extracting $k$-th roots. Thus:
(11.6) Corollary. If $E$ is a semi-ample line bundle, then $E \geqslant 0$.

In the case of vector bundles $(r \geqslant 2)$ the answer is unknown, mainly because there is no known procedure to get a Griffiths semipositive metric on $E$ from one on $S^{k} E$.

## § 11.B. Ampleness

We are now turning ourselves to the definition of ampleness. If $E \longrightarrow X$ is a holomorphic vector bundle, we define the bundle $J^{k} E$ of $k$-jets of sections of $E$ by $\left(J^{k} E\right)_{x}=\mathscr{O}_{x}(E) /\left(\mathscr{M}_{x}^{k+1} \cdot \mathscr{O}_{x}(E)\right)$ for every $x \in X$, where $\mathcal{M}_{x}$ is the maximal ideal of $\mathscr{O}_{x}$. Let $\left(e_{1}, \ldots, e_{r}\right)$ be a holomorphic frame of $E$ and $\left(z_{1}, \ldots, z_{n}\right)$ analytic coordinates
on an open subset $\Omega \subset X$. The fiber $\left(J^{k} E\right)_{x}$ can be identified with the set of Taylor developments of order $k$ :

$$
\sum_{1 \leqslant \lambda \leqslant r,|\alpha| \leqslant k} c_{\lambda, \alpha}(z-x)^{\alpha} e_{\lambda}(z),
$$

and the coefficients $c_{\lambda, \alpha}$ define coordinates along the fibers of $J^{k} E$. It is clear that the choice of another holomorphic frame ( $e_{\lambda}$ ) would yield a linear change of coordinates ( $c_{\lambda, \alpha}$ ) with holomorphic coefficients in $x$. Hence $J^{k} E$ is a holomorphic vector bundle of rank $r\binom{n+k}{n}$.

## (11.7) Definition.

a) $E$ is said to be very ample if all evaluation maps

$$
H^{0}(X, E) \rightarrow\left(J^{1} E\right)_{x}, \quad H^{0}(X, E) \rightarrow E_{x} \oplus E_{y}, \quad x, y \in X, x \neq y
$$

are surjective.
b) $E$ is said to be ample if there exists an integer $k_{0}$ such that $S^{k} E$ is very ample for $k \geqslant k_{0}$.
(11.8) Example. $\mathscr{O}(1) \rightarrow \mathbb{P}^{n}$ is a very ample line bundle (immediate verification). Since the pull-back of a (very) ample vector bundle by an embedding is clearly also (very) ample, diagram (V-16.8) shows that $\Lambda^{r} Q \rightarrow G_{r}(V)$ is very ample. However, $Q$ itself cannot be very ample if $r \geqslant 2$, because $\operatorname{dim} H^{0}\left(G_{r}(V), Q\right)=\operatorname{dim} V=d$, whereas

$$
\operatorname{rank}\left(J^{1} Q\right)=(\operatorname{rank} Q)\left(1+\operatorname{dim} G_{r}(V)\right)=r(1+r(d-r))>d \quad \text { if } r \geqslant 2
$$

(11.9) Proposition. If $E$ is very ample of rank $r$, there exists a subspace $V$ of $H^{0}(X, E)$, $\operatorname{dim} V \leqslant \max (n r+n+r, 2(n+r))$, such that all the evaluation maps $V \rightarrow E_{x} \oplus E_{y}$, $x \neq y$, and $V \rightarrow\left(J^{1} E\right)_{x}, x \in X$, are surjective.

Proof. The arguments are exactly the same as in the proof of Prop. 11.4, if we consider instead the bundles $J^{1} E \longrightarrow X$ and $E \times E \longrightarrow X \times X \backslash \Delta_{X}$ of respective ranks $r(n+1)$ and $2 r$, and sections $s_{1}^{\prime}, \ldots, s_{N}^{\prime} \in H^{0}(X, E)$ generating these bundles.
(11.10) Proposition. Let $E \rightarrow X$ be a holomorphic vector bundle.
a) If $V \subset H^{0}(X, E)$ generates $J^{1} E \longrightarrow X$ and $E \times E \longrightarrow X \times X \backslash \Delta_{X}$, then $\psi_{V}$ is an embedding.
b) Conversely, if rank $E=1$ and if there exists $V \subset H^{0}(X, E)$ generating $E$ such that $\psi_{V}$ is an embedding, then $E$ is very ample.

Proof. b) is immediate, because $E=\psi_{V}^{\star}(\mathscr{O}(1))$ and $\mathscr{O}(1)$ is very ample. Note that the result is false for $r \geqslant 2$ as shown by the example $E=Q$ over $X=G_{r}(V)$.
a) Under the assumption of a), it is clear since $S_{x}=\{s \in V ; s(x)=0\}$ that $S_{x}=S_{y}$ implies $x=y$, hence $\psi_{V}$ is injective. Therefore, it is enough to prove that the map $x \mapsto S_{x}$
has an injective differential. Let $x \in X$ and $W \subset V$ such that $S_{x} \oplus W=V$. Choose a coordinate system in a neighborhood of $x$ in $X$ and a small tangent vector $h \in T_{x} X$. The element $S_{x+h} \in G_{r}(V)$ is the graph of a small linear map $u=O(|h|): S_{x} \rightarrow W$. Thus we have

$$
S_{x+h}=\left\{s^{\prime}=s+t \in V ; s \in S_{x}, t=u(s) \in W, s^{\prime}(x+h)=0\right\}
$$

Since $s(x)=0$ and $|t|=O(|h|)$, we find

$$
s^{\prime}(x+h)=s^{\prime}(x)+d_{x} s^{\prime} \cdot h+O\left(\left|s^{\prime}\right| \cdot|h|^{2}\right)=t(x)+d_{x} s \cdot h+O\left(|s| \cdot|h|^{2}\right)
$$

thus $s^{\prime}(x+h)=0$ if and only if $t(x)=-d_{x} s \cdot h+O\left(|s| \cdot|h|^{2}\right)$. Thanks to the fiber isomorphism $\Psi_{V}: E_{x} \rightarrow V / S_{x} \simeq W, t(x) \longmapsto t \bmod S_{x}$, we get:

$$
u(s)=t=\Psi_{V}(t(x))=-\Psi_{V}\left(d_{x} s \cdot h+O\left(|s| \cdot|h|^{2}\right)\right) .
$$

Recall that $T_{y} G_{r}(V)=\operatorname{hom}\left(S_{y}, Q_{y}\right)=\operatorname{hom}(y, V / y)$ (see V-16.5) and use these identifications at $y=S_{x}$. It follows that

$$
\begin{equation*}
\left(d_{x} \psi_{V}\right) \cdot h=u=\left(S_{x} \rightarrow V / S_{x}, s \longmapsto-\Psi_{V}\left(d_{x} s \cdot h\right)\right) \tag{11.11}
\end{equation*}
$$

Now hypothesis a) implies that $S_{x} \ni s \longmapsto d_{x} s \in \operatorname{hom}\left(T_{x} X, E_{x}\right)$ is onto, hence $d_{x} \psi_{V}$ is injective.
(11.12) Corollary. If $E$ is an ample line bundle, then $E>0$.

Proof. If $E$ is very ample, diagram (11.4 ) shows that $E$ is the pull-back of $\mathscr{O}(1)$ by the embedding $\psi_{V}$, hence $i \Theta(E)=\psi_{V}^{\star}(\mathrm{i} \Theta(\mathscr{G}(1)))>0$ with the induced metric. The ample case follows by extracting roots.
(11.13) Corollary. If $E$ is a very ample vector bundle, then $E$ carries a hermitian metric such that $E^{\star}<_{\text {Nak }} 0$, in particular $E>_{\text {Grif }} 0$.

Proof. Choose $V$ as in Prop. 11.9 and select an arbitrary hermitian metric on $V$. Then diagram 11.4 yields $E=\psi_{V}^{\star} Q$, hence $\theta_{E}=\Psi_{V}^{\star} \theta_{Q}$. By formula (V-16.9) we have for every $\xi \in T G_{r}(V)=\operatorname{hom}(S, Q)$ and $t \in Q:$

$$
\theta_{Q}(\xi \otimes t, \xi \otimes t)=\sum_{j, k, l} \xi_{j k} \bar{\xi}_{l k} t_{l} \bar{t}_{j}=\sum_{k}\left|\sum_{j} \bar{t}_{j} \xi_{j k}\right|^{2}=|\langle\bullet, t\rangle \circ \xi|^{2}
$$

Let $h \in T_{x} X, t \in E_{x}$. Thanks to formula (11.11), we get

$$
\begin{aligned}
\theta_{E}(h \otimes t & h \otimes t)=\theta_{Q}\left(\left(d_{x} \psi_{V} \cdot h\right) \otimes \Psi_{V}(t),\left(d_{x} \psi_{V} \cdot h\right) \otimes \Psi_{V}(t)\right) \\
& =\left|\left\langle\bullet, \Psi_{V}(t)\right\rangle \circ\left(d_{x} \psi_{V} \cdot h\right)\right|^{2}=\left|S_{x} \ni s \longmapsto\left\langle\Psi_{V}\left(d_{x} s \cdot h\right), \Psi_{V}(t)\right\rangle\right|^{2} \\
& =\left|S_{x} \ni s \longmapsto\left\langle d_{x} s \cdot h, t\right\rangle\right|^{2} \geqslant 0 .
\end{aligned}
$$

As $S_{x} \ni s \mapsto d_{x} s \in T^{\star} X \otimes E$ is surjective, it follows that $\theta_{E}(h \otimes t, h \otimes t) \neq 0$ when $h \neq 0, t \neq 0$. Now, $d_{x} s$ defines a linear form on $T X \otimes E^{\star}$ and the above formula for the curvature of $E$ clearly yields

$$
\theta_{E^{\star}}(u, u)=-\left|S_{x} \ni s \longmapsto d_{x} s \cdot u\right|^{2}<0 \quad \text { if } u \neq 0
$$

(11.14) Problem ([Griffiths 1969]). If $E$ is an ample vector bundle over a compact manifold $X$, then is $E>_{\text {Grif }} 0$ ?

Griffiths' problem has been solved in the affirmative when $X$ is a curve (see [Umemura 1973] and also [Campana-Flenner 1990]), but the general case is still unclear and seems very deep. The next sections will be concerned with the important result of Kodaira asserting the equivalence between positivity and ampleness for line bundles.

## § 12. Blowing-up along a Submanifold

Here we generalize the blowing-up process already considered in Remark 4.5 to arbitrary manifolds. Let $X$ be a complex $n$-dimensional manifold and $Y$ a closed submanifold with $\operatorname{codim}_{X} Y=s$.
(12.1) Notations. The normal bundle of $Y$ in $X$ is the vector bundle over $Y$ defined as the quotient $N Y=(T X)_{\uparrow Y} / T Y$. The fibers of $N Y$ are thus given by $N_{y} Y=T_{y} X / T_{y} Y$ at every point $y \in Y$. We also consider the projectivized normal bundle $P(N Y) \rightarrow Y$ whose fibers are the projective spaces $P\left(N_{y} Y\right)$ associated to the fibers of $N Y$.

The blow-up of $X$ with center $Y$ (to be constructed later) is a complex $n$-dimensional manifold $\widetilde{X}$ together with a holomorphic map $\sigma: \widetilde{X} \longrightarrow X$ such that:
i) $E:=\sigma^{-1}(Y)$ is a smooth hypersurface in $\widetilde{X}$, and the restriction $\sigma: E \rightarrow Y$ is a holomorphic fiber bundle isomorphic to the projectivized normal bundle $P(N Y) \rightarrow Y$.
ii) $\sigma: \widetilde{X} \backslash E \longrightarrow X \backslash Y$ is a biholomorphism.

In order to construct $\widetilde{X}$ and $\sigma$, we first define the set-theoretic underlying objects as the disjoint sums

$$
\begin{array}{rlrl}
\widetilde{X} & =(X \backslash Y) \amalg E, & \text { where } E:=P(N Y), \\
\sigma & =\operatorname{Id}_{X \backslash Y} \amalg \pi, & & \text { where } \pi: E \longrightarrow Y .
\end{array}
$$



VII-1 Blow-up of one point in $X$.

This means intuitively that we have replaced each point $y \in Y$ by the projective space of all directions normal to $Y$. When $Y$ is reduced to a single point, the geometric picture is given by Fig. 1 below. In general, the picture is obtained by slicing $X$ transversally to $Y$ near each point and by blowing-up each slice at the intersection point with $Y$.

It remains to construct the manifold structure on $\widetilde{X}$ and in particular to describe what are the holomorphic functions near a point of $E$. Let $f, g$ be holomorphic functions on an open set $U \subset X$ such that $f=g=0$ on $Y \cap U$. Then $d f$ and $d g$ vanish on $T Y_{\upharpoonright Y \cap U}$, hence $d f$ and $d g$ induce linear forms on $N Y_{\uparrow Y \cap U}$. The holomorphic function $h(z)=f(z) / g(z)$ on the open set

$$
U_{g}:=\{z \in U ; g(z) \neq 0\} \subset U \backslash Y
$$

can be extended in a natural way to a function $\widetilde{h}$ on the set

$$
\widetilde{U}_{g}=U_{g} \cup\left\{(z,[\xi]) \in P(N Y)_{\mid Y \cap U} ; d g_{z}(\xi) \neq 0\right\} \subset \widetilde{X}
$$

by letting

$$
\widetilde{h}(z,[\xi])=\frac{d f_{z}(\xi)}{d g_{z}(\xi)}, \quad(z,[\xi]) \in P(N Y)_{\mid Y \cap U}
$$

Using this observation, we now define the manifold structure on $\widetilde{X}$ by giving explicitly an atlas. Every coordinate chart of $X \backslash Y$ is taken to be also a coordinate chart of $\widetilde{X}$. Furthermore, for every point $y_{0} \in Y$, there exists a neighborhood $U$ of $y_{0}$ in $X$ and a coordinate chart $\tau(z)=\left(z_{1}, \ldots, z_{n}\right): U \rightarrow \mathbb{C}^{n}$ centered at $y_{0}$ such that $\tau(U)=B^{\prime} \times B^{\prime \prime}$ for some balls $B^{\prime} \subset \mathbb{C}^{s}, B^{\prime \prime} \subset \mathbb{C}^{n-s}$, and such that $Y \cap U=\tau^{-1}\left(\{0\} \times B^{\prime \prime}\right)=\left\{z_{1}=\ldots=z_{s}=0\right\}$.

It follows that $\left(z_{s+1}, \ldots, z_{n}\right)$ are local coordinates on $Y \cap U$ and that the vector fields $\left(\partial / \partial z_{1}, \ldots, \partial / \partial z_{s}\right)$ yield a holomorphic frame of $N Y_{\mid Y \cap U}$. Let us denote by $\left(\xi_{1}, \ldots, \xi_{s}\right)$ the corresponding coordinates along the fibers of $N Y$. Then $\left(\xi_{1}, \ldots, \xi_{s}, z_{s+1}, \ldots, z_{n}\right)$ are coordinates on the total space $N Y$. For every $j=1, \ldots, s$, we set

$$
\widetilde{U}_{j}=\widetilde{U}_{z_{j}}=\left\{z \in U \backslash Y ; z_{j} \neq 0\right\} \cup\left\{(z,[\xi]) \in P(N Y)_{\mid Y \cap U} ; \xi_{j} \neq 0\right\} .
$$

Then $\left(\widetilde{U}_{j}\right)_{1 \leqslant j \leqslant s}$ is a covering of $\widetilde{U}=\sigma^{-1}(U)$ and for each $j$ we define a coordinate chart $\widetilde{\tau}_{j}=\left(w_{1}, \ldots, w_{n}\right): \widetilde{U}_{j} \longrightarrow \mathbb{C}^{n}$ by

$$
w_{k}:=\left(\frac{z_{k}}{z_{j}}\right)^{\sim} \quad \text { for } 1 \leqslant k \leqslant s, k \neq j ; \quad w_{k}:=z_{k} \quad \text { for } k>s \text { or } k=j .
$$

For $z \in U \backslash Y$, resp. $(z,[\xi]) \in P(N Y)_{\mid Y \cap U}$, we get

$$
\begin{aligned}
& \widetilde{\tau}_{j}(z)=\left(w_{1}, \ldots, w_{n}\right)=\left(\frac{z_{1}}{z_{j}}, \ldots, \frac{z_{j-1}}{z_{j}}, z_{j}, \frac{z_{j+1}}{z_{j}}, \ldots, \frac{z_{s}}{z_{j}}, z_{s+1}, \ldots, z_{n}\right), \\
& \widetilde{\tau}_{j}(z,[\xi])=\left(w_{1}, \ldots, w_{n}\right)=\left(\frac{\xi_{1}}{\xi_{j}}, \ldots, \frac{\xi_{j-1}}{\xi_{j}}, 0, \frac{\xi_{j+1}}{\xi_{j}}, \ldots, \frac{\xi_{s}}{\xi_{j}}, \xi_{s+1}, \ldots, \xi_{n}\right) .
\end{aligned}
$$

With respect to the coordinates $\left(w_{k}\right)$ on $\widetilde{U}_{j}$ and $\left(z_{k}\right)$ on $U$, the map $\sigma$ is given by

$$
\begin{align*}
\widetilde{U}_{j} & \stackrel{\sigma}{\longrightarrow} U  \tag{12.2}\\
w & \stackrel{\sigma_{j}}{\longmapsto}\left(w_{1} w_{j}, \ldots, w_{j-1} w_{j} ; w_{j} ; w_{j+1} w_{j}, \ldots, w_{s} w_{j} ; w_{s+1}, \ldots, w_{n}\right)
\end{align*}
$$

where $\sigma_{j}=\tau \circ \sigma \circ \widetilde{\tau}_{j}^{-1}$, thus $\sigma$ is holomorphic. The range of the coordinate chart $\widetilde{\tau}_{j}$ is $\widetilde{\tau}_{j}\left(\widetilde{U}_{j}\right)=\sigma_{j}^{-1}(\tau(U))$, so it is actually open in $\mathbb{C}^{n}$. Furthermore $E \cap \widetilde{U}_{j}$ is defined by the single equation $w_{j}=0$, thus $E$ is a smooth hypersurface in $\widetilde{X}$. It remains only to verify that the coordinate changes $w \longmapsto w^{\prime}$ associated to any coordinate change $z \longmapsto z^{\prime}$ on $X$ are holomorphic. For that purpose, it is sufficient to verify that $(f / g)^{\sim}$ is holomorphic in $\left(w_{1}, \ldots, w_{n}\right)$ on $\widetilde{U}_{j} \cap \widetilde{U}_{g}$. As $g$ vanishes on $Y \cap U$, we can write $g(z)=\sum_{1 \leqslant k \leqslant s} z_{k} A_{k}(z)$ for some holomorphic functions $A_{k}$ on $U$. Therefore

$$
\frac{g(z)}{z_{j}}=A_{j}\left(\sigma_{j}(w)\right)+\sum_{k \neq j} w_{k} A_{k}\left(\sigma_{j}(w)\right)
$$

has an extension $\left(g / z_{j}\right)^{\sim}$ to $\widetilde{U}_{j}$ which is a holomorphic function of the variables $\left(w_{1}, \ldots, w_{n}\right)$. Since $\left(g / z_{j}\right)^{\sim}(z,[\xi])=d g_{z}(\xi) / \xi_{j}$ on $E \cap \widetilde{U}_{j}$, it is clear that

$$
\widetilde{U}_{j} \cap \widetilde{U}_{g}=\left\{w \in \widetilde{U}_{j} ;\left(g / z_{j}\right)^{\sim}(w) \neq 0\right\} .
$$

Hence $\widetilde{U}_{j} \cap \widetilde{U}_{g}$ is open in $\widetilde{U}_{g}$ and $(f / g)^{\sim}=\left(f / z_{j}\right)^{\sim} /\left(g / z_{j}\right)^{\sim}$ is holomorphic on $\widetilde{U}_{j} \cap \widetilde{U}_{g}$.
(12.3) Definition. The map $\sigma: \widetilde{X} \rightarrow X$ is called the blow-up of $X$ with center $Y$ and $E=\sigma^{-1}(Y) \simeq P(N Y)$ is called the exceptional divisor of $\widetilde{X}$.

According to (V-13.5), we denote by $\mathscr{O}(E)$ the line bundle on $\widetilde{X}$ associated to the divisor $E$ and by $h \in H^{0}(\widetilde{X}, \mathscr{O}(E))$ the canonical section such that $\operatorname{div}(h)=[E]$. On the other hand, we denote by $\mathscr{O}_{P(N Y)}(-1) \subset \pi^{\star}(N Y)$ the tautological line subbundle over $E=P(N Y)$ such that the fiber above the point $(z,[\xi])$ is $\mathbb{C} \xi \subset N_{z} Y$.
(12.4) Proposition. $\mathfrak{G}(E)$ enjoys the following properties:
a) $\mathscr{O}(E)_{\upharpoonright E}$ is isomorphic to $\mathscr{O}_{P(N Y)}(-1)$.
b) Assume that $X$ is compact. For every positive line bundle $L$ over $X$, the line bundle $\mathfrak{O}(-E) \otimes \sigma^{\star}\left(L^{k}\right)$ over $\widetilde{X}$ is positive for $k>0$ large enough.

Proof. a) The canonical section $h \in H^{0}(\widetilde{X}, \mathscr{O}(E))$ vanishes at order 1 along $E$, hence the kernel of its differential

$$
d h:(T \widetilde{X})_{\mid E} \longrightarrow \mathscr{O}(E)_{\mid E}
$$

is $T E$. We get therefore an isomorphism $N E \simeq \mathscr{O}(E)_{\mid E}$. Now, the map $\sigma: \widetilde{X} \rightarrow X$ satisfies $\sigma(E) \subset Y$, so its differential $d \sigma: T \widetilde{X} \longrightarrow \sigma^{\star}(T X)$ is such that $d \sigma(T E) \subset$ $\sigma^{\star}(T Y)$. Therefore $d \sigma$ induces a morphism

$$
\begin{equation*}
N E \longrightarrow \sigma^{\star}(N Y)=\pi^{\star}(N Y) \tag{12.5}
\end{equation*}
$$

of vector bundles over $E$. The vector field $\partial / \partial w_{j}$ yields a non vanishing section of $N E$ on $\widetilde{U}_{j}$, and (12.2) implies

$$
d \sigma_{j}\left(\frac{\partial}{\partial w_{j}}\right)=\frac{\partial}{\partial z_{j}}+\sum_{1 \leqslant k \leqslant s, k \neq j} w_{k} \frac{\partial}{\partial z_{k}} \quad / / \sum_{1 \leqslant k \leqslant s} \xi_{k} \frac{\partial}{\partial z_{k}}
$$

at every point $(z,[\xi]) \in E$. This shows that (12.5) is an isomorphism of $N E$ onto $\mathscr{O}_{P(N Y)}(-1) \subset \pi^{\star}(N Y)$, hence

$$
\begin{equation*}
\mathscr{O}(E)_{\upharpoonright E} \simeq N E \simeq \mathscr{O}_{P(N Y)}(-1) . \tag{12.6}
\end{equation*}
$$

b) Select an arbitrary hermitian metric on $T X$ and consider the induced metrics on $N Y$ and on $\mathscr{O}_{P(N Y)}(1) \longrightarrow E=P(N Y)$. The restriction of $\mathscr{O}_{P(N Y)}(1)$ to each fiber $P\left(N_{z} Y\right)$ is the standard line bundle $\mathscr{O}(1)$ over $\mathbb{P}^{s-1}$; thus by (V-15.10) this restriction has a positive definite curvature form. Extend now the metric of $\mathscr{O}_{P(N Y)}(1)$ on $E$ to a metric of $\mathbb{O}(-E)$ on $X$ in an arbitrary way. If $F=\mathscr{O}(-E) \otimes \sigma^{\star}\left(L^{k}\right)$, then $\Theta(F)=\Theta(\mathscr{O}(-E))+k \sigma^{\star} \Theta(L)$, thus for every $t \in T \widetilde{X}$ we have

$$
\theta_{F}(t, t)=\theta_{\sigma(-E)}(t, t)+k \theta_{L}(d \sigma(t), d \sigma(t))
$$

By the compactness of the unitary tangent bundle to $\widetilde{X}$ and the positivity of $\theta_{L}$, it is sufficient to verify that $\theta_{\widehat{O}(-E)}(t, t)>0$ for every unit vector $t \in T_{z} \widetilde{X}$ such that $d \sigma(t)=0$. However, from the computations of a), this can only happen when $z \in E$ and $t \in T E$, and in that case $d \sigma(t)=d \pi(t)=0$, so $t$ is tangent to the fiber $P\left(N_{z} Y\right)$. Therefore

$$
\theta_{\overparen{O}(-E)}(t, t)=\theta_{\widehat{G}_{P(N Y)}(1)}(t, t)>0
$$

(12.7) Proposition. The canonical line bundle of $\widetilde{X}$ is given by

$$
K_{\widetilde{X}}=\overparen{O}((s-1) E) \otimes \sigma^{\star} K_{X}, \quad \text { where } \quad s=\operatorname{codim}_{X} Y .
$$

Proof. $K_{X}$ is generated on $U$ by the holomorphic $n$-form $d z_{1} \wedge \ldots \wedge d z_{n}$. Using (12.2), we see that $\sigma^{\star} K_{X}$ is generated on $\widetilde{U}_{j}$ by

$$
\sigma^{\star}\left(d z_{1} \wedge \ldots \wedge d z_{n}\right)=w_{j}^{s-1} d w_{1} \wedge \ldots \wedge d w_{n}
$$

Since the divisor of the section $h \in H^{0}(\widetilde{X}, \mathscr{O}(E))$ is the hypersurface $E$ defined by the equation $w_{j}=0$ in $\widetilde{U}_{j}$, we have a well defined line bundle isomorphism

$$
\sigma^{\star} K_{X} \longrightarrow \mathscr{O}((1-s) E) \otimes K_{\widetilde{X}}, \quad \alpha \longmapsto h^{1-s} \sigma^{\star}(\alpha)
$$

## § 13. Equivalence of Positivity and Ampleness for Line Bundles

We have seen in section 11 that every ample line bundle carries a hermitian metric of positive curvature. The converse will be a consequence of the following result.
(13.1) Theorem. Let $L \longrightarrow X$ be a positive line bundle and $L^{k}$ the $k$-th tensor power of $L$. For every $N$-tuple $\left(x_{1}, \ldots, x_{N}\right)$ of distinct points of $X$, there exists a constant $C>0$ such that the evaluation maps

$$
H^{0}\left(X, L^{k}\right) \longrightarrow\left(J^{m} L^{k}\right)_{x_{1}} \oplus \cdots \oplus\left(J^{m} L^{k}\right)_{x_{N}}
$$

are surjective for all integers $m \geqslant 0, k \geqslant C(m+1)$.
(13.2) Lemma. Let $\sigma: \widetilde{X} \longrightarrow X$ be the blow-up of $X$ with center the finite set $Y=$ $\left\{x_{1}, \ldots, x_{N}\right\}$, and let $\mathfrak{O}(E)$ be the line bundle associated to the exceptional divisor $E$. Then

$$
H^{1}\left(\widetilde{X}, \overparen{O}(-m E) \otimes \sigma^{\star} L^{k}\right)=0
$$

for $m \geqslant 1, k \geqslant C m$ and $C \geqslant 0$ large enough.
Proof. By Prop. 12.7 we get $K_{\widetilde{X}}=\mathscr{O}((n-1) E) \otimes \sigma^{\star} K_{X}$ and

$$
H^{1}\left(\widetilde{X}, \mathscr{O}(-m E) \otimes \sigma^{\star} L^{k}\right)=H^{n, 1}\left(\widetilde{X}, K_{\widetilde{X}}^{-1} \otimes \mathbb{O}(-m E) \otimes \sigma^{\star} L^{k}\right)=H^{n, 1}(\widetilde{X}, F)
$$

where $F=\mathscr{O}(-(m+n-1) E) \otimes \sigma^{\star}\left(K_{X}^{-1} \otimes L^{k}\right)$, so the conclusion will follow from the Kodaira-Nakano vanishing theorem if we can show that $F>0$ when $k$ is large enough. Fix an arbitrary hermitian metric on $K_{X}$. Then

$$
\Theta(F)=(m+n-1) \Theta(\mathbb{O}(-E))+\sigma^{\star}\left(k \Theta(L)-\Theta\left(K_{X}\right)\right) .
$$

There is $k_{0} \geqslant 0$ such that $\mathrm{i}\left(k_{0} \Theta(L)-\Theta\left(K_{X}\right)\right)>0$ on $X$, and Prop. 12.4 implies the existence of $C_{0}>0$ such that $\mathrm{i}\left(\Theta(\mathscr{O}(-E))+C_{0} \sigma^{\star} \Theta(L)\right)>0$ on $\widetilde{X}$. Thus $\mathrm{i} \Theta(F)>0$ for $m \geqslant 2-n$ and $k \geqslant k_{0}+C_{0}(m+n-1)$.

Proof of Theorem 13.1.. Let $v_{j} \in H^{0}\left(\Omega_{j}, L^{k}\right)$ be a holomorphic section of $L^{k}$ in a neighborhood $\Omega_{j}$ of $x_{j}$ having a prescribed $m$-jet at $x_{j}$. Set

$$
v(x)=\sum_{j} \psi_{j}(x) v_{j}(x)
$$

where $\psi_{j}=1$ in a neighborhood of $x_{j}$ and $\psi_{j}$ has compact support in $\Omega_{j}$. Then $d^{\prime \prime} v=$ $\sum d^{\prime \prime} \psi_{j} \cdot v_{j}$ vanishes in a neighborhood of $x_{1}, \ldots, x_{N}$. Let $h$ be the canonical section of $\mathscr{O}(E)^{-1}$ such that $\operatorname{div}(h)=[E]$. The $(0,1)$-form $\sigma^{\star} d^{\prime \prime} v$ vanishes in a neighborhood of $E=h^{-1}(0)$, hence

$$
w=h^{-(m+1)} \sigma^{\star} d^{\prime \prime} v \in \mathscr{C}_{0,1}^{\infty}\left(\widetilde{X}, \overparen{O}(-(m+1) E) \otimes \sigma^{\star} L^{k}\right)
$$

and $w$ is a $d^{\prime \prime}$-closed form. By Lemma 13.2 there exists a smooth section

$$
u \in \mathscr{C}_{0,0}^{\infty}\left(\widetilde{X}, \mathscr{O}(-(m+1) E) \otimes \sigma^{\star} L^{k}\right)
$$

such that $d^{\prime \prime} u=w=h^{-(m+1)} \sigma^{\star} d^{\prime \prime} v$. This implies

$$
\sigma^{\star} v-h^{m+1} u \in H^{0}\left(\widetilde{X}, \sigma^{\star} L^{k}\right)
$$

and since $\sigma^{\star} L$ is trivial near $E$, there exists a section $g \in H^{0}\left(X, L^{k}\right)$ such that $\sigma^{\star} g=$ $\sigma^{\star} v-h^{m+1} u$. As $h$ vanishes at order 1 along $E$, the $m$-jet of $g$ at $x_{j}$ must be equal to that of $v$ (or $v_{j}$ ).
(13.3) Corollary. For any holomorphic line bundle $L \longrightarrow X$, the following conditions are equivalent:
a) $L$ is ample;
b) $L>0$, i.e. $L$ possesses a hermitian metric such that $\mathrm{i} \Theta(L)>0$.

Proof. a) $\Longrightarrow \mathrm{b}$ ) is given by Cor. 11.12, whereas b$) \Longrightarrow \mathrm{a}$ ) is a consequence of Th. 13.1 for $m=1$.

## § 14. Kodaira's Projectivity Criterion

The following fundamental projectivity criterion is due to [Kodaira 1954].
(14.1) Theorem. Let $X$ be a compact complex manifold, $\operatorname{dim}_{\mathbb{C}} X=n$. The following conditions are equivalent.
a) $X$ is projective algebraic, i.e. $X$ can be embedded as an algebraic submanifold of the complex projective space $\mathbb{P}^{N}$ for $N$ large.
b) $X$ carries a positive line bundle $L$.
c) $X$ carries a Hodge metric, i.e. a Kähler metric $\omega$ with rational cohomology class $\{\omega\} \in H^{2}(X, \mathbb{Q})$.

Proof. a) $\Longrightarrow \mathrm{b})$. Take $L=\mathscr{O}(1)_{\mid X}$.
$\mathrm{b}) \Longrightarrow \mathrm{c})$. Take $\omega=\frac{\mathrm{i}}{2 \pi} \Theta(L)$; then $\{\omega\}$ is the image of $c_{1}(L) \in H^{2}(X, \mathbb{Z})$.
c) $\Longrightarrow$ b). We can multiply $\{\omega\}$ by a common denominator of its coefficients and suppose that $\{\omega\}$ is in the image of $H^{2}(X, \mathbb{Z})$. Then Th. V-13.9 b) shows that there exists a hermitian line bundle $L$ such that $\frac{i}{2 \pi} \Theta(L)=\omega>0$.
b) $\Longrightarrow$ a). Corollary 13.3 shows that $F=L^{k}$ is very ample for some integer $k>0$. Then Prop. 11.9 enables us to find a subspace $V$ of $H^{0}(X, F)$, $\operatorname{dim} V \leqslant 2 n+2$, such that $\psi_{V}: X \longrightarrow G_{1}(V)=P\left(V^{\star}\right)$ is an embedding. Thus $X$ can be embedded in $\mathbb{P}^{2 n+1}$ and Chow's theorem II-7.10 shows that the image is an algebraic set in $\mathbb{P}^{2 n+1}$.
(14.2) Remark. The above proof shows in particular that every $n$-dimensional projective manifold $X$ can be embedded in $\mathbb{P}^{2 n+1}$. This can be shown directly by using generic projections $\mathbb{P}^{N} \rightarrow \mathbb{P}^{2 n+1}$ and Whitney type arguments as in 11.2.
(14.3) Corollary. Every compact Riemann surface $X$ is isomorphic to an algebraic curve in $\mathbb{P}^{3}$.

Proof. Any positive smooth form $\omega$ of type $(1,1)$ is Kähler, and $\omega$ is in fact a Hodge metric if we normalize its volume so that $\int_{X} \omega=1$.

This example can be somewhat generalized as follows.
(14.4) Corollary. Every Kähler manifold $(X, \omega)$ such that $H^{2}(X, \mathscr{O})=0$ is projective.

Proof. By hypothesis $H^{0,2}(X, \mathbb{C})=0=H^{2,0}(X, \mathbb{C})$, hence

$$
H^{2}(X, \mathbb{C})=H^{1,1}(X, \mathbb{C})
$$

admits a basis $\left\{\alpha_{1}\right\}, \ldots,\left\{\alpha_{N}\right\} \in H^{2}(X, \mathbb{Q})$ where $\alpha_{1}, \ldots, \alpha_{N}$ are harmonic real $(1,1)$ forms. Since $\{\omega\}$ is real, we have $\{\omega\}=\lambda_{1}\left\{\alpha_{1}\right\}+\ldots+\lambda_{N}\left\{\alpha_{N}\right\}, \lambda_{j} \in \mathbb{R}$, thus

$$
\omega=\lambda_{1} \alpha_{1}+\ldots+\lambda_{N} \alpha_{N}
$$

because $\omega$ itself is harmonic. If $\mu_{1}, \ldots, \mu_{N}$ are rational numbers sufficiently close to $\lambda_{1}, \ldots, \lambda_{N}$, then $\widetilde{\omega}:=\mu_{1} \alpha_{1}+\cdots \mu_{N} \alpha_{N}$ is close to $\omega$, so $\widetilde{\omega}$ is a positive definite $d$-closed (1, 1)-form, and $\{\widetilde{\omega}\} \in H^{2}(X, \mathbb{Q})$.

We obtain now as a consequence the celebrated Riemann criterion characterizing abelian varieties ( $=$ projective algebraic complex tori).
(14.5) Corollary. A complex torus $X=\mathbb{C}^{n} / \Gamma\left(\Gamma\right.$ a lattice of $\left.\mathbb{C}^{n}\right)$ is an abelian variety if and only if there exists a positive definite hermitian form $h$ on $\mathbb{C}^{n}$ such that

$$
\operatorname{Im}\left(h\left(\gamma_{1}, \gamma_{2}\right)\right) \in \mathbb{Z} \quad \text { for all } \gamma_{1}, \gamma_{2} \in \Gamma
$$

Proof (Sufficiency of the condition).. Set $\omega=-\operatorname{Im} h$. Then $\omega$ defines a constant Kähler metric on $\mathbb{C}^{n}$, hence also on $X=\mathbb{C}^{n} / \Gamma$. Let $\left(a_{1}, \ldots, a_{2 n}\right)$ be an integral basis of the lattice $\Gamma$. We denote by $T_{j}, T_{j k}$ the real 1- and 2-tori

$$
T_{j}=(\mathbb{R} / \mathbb{Z}) a_{j}, \quad 1 \leqslant j \leqslant n, \quad T_{j k}=T_{j} \oplus T_{k}, \quad 1 \leqslant j<k \leqslant 2 n .
$$

Topologically we have $X \approx T_{1} \times \ldots \times T_{2 n}$, so the Künneth formula IV-15.7 yields

$$
\begin{aligned}
& H^{\bullet}(X, \mathbb{Z}) \simeq \bigotimes_{1 \leqslant j \leqslant 2 n}\left(H^{0}\left(T_{j}, \mathbb{Z}\right) \oplus H^{1}\left(T_{j}, \mathbb{Z}\right)\right), \\
& H^{2}(X, \mathbb{Z}) \simeq \bigoplus_{1 \leqslant j<k \leqslant 2 n} H^{1}\left(T_{j}, \mathbb{Z}\right) \otimes H^{1}\left(T_{k}, \mathbb{Z}\right) \simeq \bigoplus_{1 \leqslant j<k \leqslant 2 n} H^{2}\left(T_{j k}, \mathbb{Z}\right)
\end{aligned}
$$

where the projection $H^{2}(X, \mathbb{Z}) \longrightarrow H^{2}\left(T_{j k}, \mathbb{Z}\right)$ is induced by the injection $T_{j k} \subset X$. In the identification $H^{2}\left(T_{j k}, \mathbb{R}\right) \simeq \mathbb{R}$, we get

$$
\begin{equation*}
\{\omega\}_{\upharpoonright T_{j k}}=\int_{T_{j k}} \omega=\omega\left(a_{j}, a_{k}\right)=-\operatorname{Im} h\left(a_{j}, a_{k}\right) . \tag{14.6}
\end{equation*}
$$

The assumption on $h$ implies $\{\omega\}_{\mid T_{j k}} \in H^{2}\left(T_{j k}, \mathbb{Z}\right)$ for all $j, k$, therefore $\{\omega\} \in H^{2}(X, \mathbb{Z})$ and $X$ is projective by Th. (14.1).

Proof (Necessity of the condition).. If $X$ is projective, then $X$ admits a Kähler metric $\omega$ such that $\{\omega\}$ is in the image of $H^{2}(X, \mathbb{Z})$. In general, $\omega$ is not invariant under the translations $\tau_{x}(y)=y-x$ of $X$. Therefore, we replace $\omega$ by its "mean value":

$$
\widetilde{\omega}=\frac{1}{\operatorname{Vol}(X)} \int_{x \in X}\left(\tau_{x}^{\star} \omega\right) d x
$$

which has the same cohomology class as $\omega$ ( $\tau_{x}$ is homotopic to the identity). Now $\widetilde{\omega}$ is the imaginary part of a constant positive definite hermitian form $h$ on $\mathbb{C}^{n}$, and formula (14.6) shows that $\operatorname{Im} h\left(a_{j}, a_{k}\right) \in \mathbb{Z}$.
(14.7) Example. Let $X$ be a projective manifold. We shall prove that the Jacobian $\operatorname{Jac}(X)$ and the Albanese variety $\operatorname{Alb}(X)$ (cf. § VI-13 for definitions) are abelian varieties.

In fact, let $\omega$ be a Kähler metric on $X$ such that $\{\omega\}$ is in the image of $H^{2}(X, \mathbb{Z})$ and let $h$ be the hermitian metric on $H^{1}(X, \mathscr{C}) \simeq H^{0,1}(X, \mathbb{C})$ defined by

$$
h(u, v)=\int_{X}-2 \mathrm{i} u \wedge \bar{v} \wedge \omega^{n-1}
$$

for all closed ( 0,1 )-forms $u, v$. As

$$
-2 \mathrm{i} u \wedge \bar{v} \wedge \omega^{n-1}=\frac{2}{n}|u|^{2} \omega^{n}
$$

we see that $h$ is a positive definite hermitian form on $H^{0,1}(X, \mathbb{C})$. Consider elements $\gamma_{j} \in H^{1}(X, \mathbb{Z}), j=1,2$. If we write $\gamma_{j}=\gamma_{j}^{\prime}+\gamma_{j}^{\prime \prime}$ in the decomposition $H^{1}(X, \mathbb{C})=$ $H^{1,0}(X, \mathbb{C}) \oplus H^{0,1}(X, \mathbb{C})$, we get

$$
\begin{aligned}
h\left(\gamma_{1}^{\prime \prime}, \gamma_{2}^{\prime \prime}\right) & =\int_{X}-2 \mathrm{i} \gamma_{1}^{\prime \prime} \wedge \gamma_{2}^{\prime} \wedge \omega^{n-1} \\
\operatorname{Im} h\left(\gamma_{1}^{\prime \prime}, \gamma_{2}^{\prime \prime}\right) & =\int_{X}\left(\gamma_{1}^{\prime} \wedge \gamma_{2}^{\prime \prime}+\gamma_{1}^{\prime \prime} \wedge \gamma_{2}^{\prime}\right) \wedge \omega^{n-1}=\int_{X} \gamma_{1} \wedge \gamma_{2} \wedge \omega^{n-1} \in \mathbb{Z}
\end{aligned}
$$

Therefore $\operatorname{Jac}(X)$ is an abelian variety.
Now, we observe that $H^{n-1, n}(X, \mathbb{C})$ is the anti-dual of $H^{0,1}(X, \mathbb{C})$ by Serre duality. We select on $H^{n-1, n}(X, \mathbb{C})$ the dual hermitian metric $h^{\star}$. Since the Poincaré bilinear pairing yields a unimodular bilinear map

$$
H^{1}(X, \mathbb{Z}) \times H^{2 n-1}(X, \mathbb{Z}) \longrightarrow \mathbb{Z}
$$

we easily conclude that $\operatorname{Im} h^{\star}\left(\gamma_{1}^{\prime \prime}, \gamma_{2}^{\prime \prime}\right) \in \mathbb{Q}$ for all $\gamma_{1}, \gamma_{2} \in H^{2 n-1}(X, \mathbb{Z})$. Therefore $\operatorname{Alb}(X)$ is also an abelian variety.

## Chapter VIII

## $L^{2}$ Estimates on Pseudoconvex Manifolds


#### Abstract

The main goal of this chapter is to show that the differential geometric technique that has been used in order to prove vanishing theorems also yields very precise $L^{2}$ estimates for the solutions of equations $d^{\prime \prime} u=v$ on pseudoconvex manifolds. The central idea, due to [Hörmander 1965], is to introduce weights of the type $e^{-\varphi}$ where $\varphi$ is a function satisfying suitable convexity conditions. This method leads to generalizations of many standard vanishing theorems to weakly pseudoconvex manifolds. As a special case, we obtain the original Hörmander estimates for pseudoconvex domains of $\mathbb{C}^{n}$, and give some applications to algebraic geometry (Hörmander-Bombieri-Skoda theorem, properties of zero sets of polynomials in $\left.\mathbb{C}^{n}\right)$. We also derive the Ohsawa-Takegoshi extension theorem for $L^{2}$ holomorphic functions and Skoda's $L^{2}$ estimates for surjective bundle morphisms [Skoda 1972a, 1978], [Demailly 1982c]. Skoda's estimates can be used to obtain a quick solution of the Levi problem, and have important applications to local algebra and Nullstellensatz theorems. Finally, $L^{2}$ estimates are used to prove the Newlander-Nirenberg theorem on the analyticity of almost complex structures. We apply it to establish Kuranishi's theorem on deformation theory of compact complex manifolds.


## § 1. Non Bounded Operators on Hilbert Spaces

A few preliminaries of functional analysis will be needed here. Let $\mathscr{H}_{1}, \mathscr{H}_{2}$ be complex Hilbert spaces. We consider a linear operator $T$ defined on a subspace Dom $T \subset \mathscr{H}_{1}$ (called the domain of $T$ ) into $\mathscr{H}_{2}$. The operator $T$ is said to be densely defined if Dom $T$ is dense in $\mathscr{H}_{1}$, and closed if its graph

$$
\operatorname{Gr} T=\{(x, T x) ; x \in \operatorname{Dom} T\}
$$

is closed in $\mathscr{H}_{1} \times \mathscr{H}_{2}$.
Assume now that $T$ is closed and densely defined. The adjoint $T^{\star}$ of $T$ (in Von Neumann's sense) is constructed as follows: Dom $T^{\star}$ is the set of $y \in \mathscr{H}_{2}$ such that the linear form

$$
\operatorname{Dom} T \ni x \longmapsto\langle T x, y\rangle_{2}
$$

is bounded in $\mathscr{H}_{1}$-norm. Since $\operatorname{Dom} T$ is dense, there exists for every $y$ in $\operatorname{Dom} T^{\star}$ a unique element $T^{\star} y \in \mathscr{H}_{1}$ such that $\langle T x, y\rangle_{2}=\left\langle x, T^{\star} y\right\rangle_{1}$ for all $x \in \operatorname{Dom} T^{\star}$. It is immediate to verify that $\operatorname{Gr} T^{\star}=(\operatorname{Gr}(-T))^{\perp}$ in $\mathscr{H}_{1} \times \mathscr{H}_{2}$. It follows that $T^{\star}$ is closed and that every pair $(u, v) \in \mathscr{H}_{1} \times \mathscr{H}_{2}$ can be written

$$
(u, v)=(x,-T x)+\left(T^{\star} y, y\right), \quad x \in \operatorname{Dom} T, \quad y \in \operatorname{Dom} T^{\star} .
$$

Take in particular $u=0$. Then

$$
x+T^{\star} y=0, \quad v=y-T x=y+T T^{\star} y, \quad\langle v, y\rangle_{2}=\|y\|_{2}^{2}+\left\|T^{\star} y\right\|_{1}^{2} .
$$

If $v \in\left(\operatorname{Dom} T^{\star}\right)^{\perp}$ we get $\langle v, y\rangle_{2}=0$, thus $y=0$ and $v=0$. Therefore $T^{\star}$ is densely defined and our discussion implies:
(1.1) Theorem [Von Neumann 1929]). If $T: \mathscr{H}_{1} \longrightarrow \mathscr{H}_{2}$ is a closed and densely defined operator, then its adjoint $T^{\star}$ is also closed and densely defined and $\left(T^{\star}\right)^{\star}=T$. Furthermore, we have the relation $\operatorname{Ker} T^{\star}=(\operatorname{Im} T)^{\perp}$ and its dual $(\operatorname{Ker} T)^{\perp}=\overline{\operatorname{Im} T^{\star}}$.

Consider now two closed and densely defined operators $T, S$ :

$$
\mathscr{H}_{1} \xrightarrow{T} \mathscr{H}_{2} \xrightarrow{S} \mathscr{H}_{3}
$$

such that $S \circ T=0$. By this, we mean that the range $T(\operatorname{Dom} T)$ is contained in $\operatorname{Ker} S \subset$ $\operatorname{Dom} S$, in such a way that there is no problem for defining the composition $S \circ T$. The starting point of all $L^{2}$ estimates is the following abstract existence theorem.
(1.2) Theorem. There are orthogonal decompositions

$$
\begin{aligned}
\mathscr{H}_{2} & =\left(\operatorname{Ker} S \cap \operatorname{Ker} T^{\star}\right) \oplus \overline{\operatorname{Im} T} \oplus \overline{\operatorname{Im} S^{\star}}, \\
\operatorname{Ker} S & =\left(\operatorname{Ker} S \cap \operatorname{Ker} T^{\star}\right) \oplus \overline{\operatorname{Im} T} .
\end{aligned}
$$

In order that $\operatorname{Im} T=\operatorname{Ker} S$, it suffices that

$$
\begin{equation*}
\left\|T^{\star} x\right\|_{1}^{2}+\|S x\|_{3}^{2} \geqslant C\|x\|_{2}^{2}, \quad \forall x \in \operatorname{Dom} S \cap \operatorname{Dom} T^{\star} \tag{1.3}
\end{equation*}
$$

for some constant $C>0$. In that case, for every $v \in \mathscr{H}_{2}$ such that $S v=0$, there exists $u \in \mathscr{H}_{1}$ such that $T u=v$ and

$$
\|u\|_{1}^{2} \leqslant \frac{1}{C}\|v\|_{2}^{2}
$$

In particular

$$
\overline{\operatorname{Im} T}=\operatorname{Im} T=\operatorname{Ker} S, \quad \overline{\operatorname{Im} S^{\star}}=\operatorname{Im} S^{\star}=\operatorname{Ker} T^{\star}
$$

Proof. Since $S$ is closed, the kernel $\operatorname{Ker} S$ is closed in $\mathscr{H}_{2}$. The relation $(\operatorname{Ker} S)^{\perp}=\overline{\operatorname{Im} S^{\star}}$ implies

$$
\begin{equation*}
\mathscr{H}_{2}=\operatorname{Ker} S \oplus \overline{\operatorname{Im} S^{\star}} \tag{1.4}
\end{equation*}
$$

and similarly $\mathscr{H}_{2}=\operatorname{Ker} T^{\star} \oplus \overline{\operatorname{Im} T}$. However, the assumption $S \circ T=0$ shows that $\overline{\operatorname{Im} T} \subset \operatorname{Ker} S$, therefore

$$
\begin{equation*}
\operatorname{Ker} S=\left(\operatorname{Ker} S \cap \operatorname{Ker} T^{\star}\right) \oplus \overline{\operatorname{Im} T} . \tag{1.5}
\end{equation*}
$$

The first two equalities in Th. 1.2 are then equivalent to the conjunction of (1.4) and (1.5).

Now, under assumption (1.3), we are going to show that the equation $T u=v$ is always solvable if $S v=0$. Let $x \in \operatorname{Dom} T^{\star}$. One can write

$$
x=x^{\prime}+x^{\prime \prime} \quad \text { where } x^{\prime} \in \operatorname{Ker} S \text { and } x^{\prime \prime} \in(\operatorname{Ker} S)^{\perp} \subset(\operatorname{Im} T)^{\perp}=\operatorname{Ker} T^{\star} .
$$

Since $x, x^{\prime \prime} \in \operatorname{Dom} T^{\star}$, we have also $x^{\prime} \in \operatorname{Dom} T^{\star}$. We get

$$
\langle v, x\rangle_{2}=\left\langle v, x^{\prime}\right\rangle_{2}+\left\langle v, x^{\prime \prime}\right\rangle_{2}=\left\langle v, x^{\prime}\right\rangle_{2}
$$

because $v \in \operatorname{Ker} S$ and $x^{\prime \prime} \in(\operatorname{Ker} S)^{\perp}$. As $S x^{\prime}=0$ and $T^{\star} x^{\prime \prime}=0$, the Cauchy-Schwarz inequality combined with (1.3) implies

$$
\left|\langle v, x\rangle_{2}\right|^{2} \leqslant\|v\|_{2}^{2}\left\|x^{\prime}\right\|_{2}^{2} \leqslant \frac{1}{C}\|v\|_{2}^{2}\left\|T^{\star} x^{\prime}\right\|_{1}^{2}=\frac{1}{C}\|v\|_{2}^{2}\left\|T^{\star} x\right\|_{1}^{2}
$$

This shows that the linear form $T_{X}^{\star} \ni x \longmapsto\langle x, v\rangle_{2}$ is continuous on $\operatorname{Im} T^{\star} \subset \mathscr{H}_{1}$ with norm $\leqslant C^{-1 / 2}\|v\|_{2}$. By the Hahn-Banach theorem, this form can be extended to a continuous linear form on $\mathscr{H}_{1}$ of norm $\leqslant C^{-1 / 2}\|v\|_{2}$, i.e. we can find $u \in \mathscr{H}_{1}$ such that $\|u\|_{1} \leqslant C^{-1 / 2}\|v\|_{2}$ and

$$
\langle x, v\rangle_{2}=\left\langle T^{\star} x, u\right\rangle_{1}, \quad \forall x \in \operatorname{Dom} T^{\star} .
$$

This means that $u \in \operatorname{Dom}\left(T^{\star}\right)^{\star}=\operatorname{Dom} T$ and $v=T u$. We have thus shown that $\operatorname{Im} T=\operatorname{Ker} S$, in particular $\operatorname{Im} T$ is closed. The dual equality $\operatorname{Im} S^{\star}=\operatorname{Ker} T^{\star}$ follows by considering the dual pair $\left(S^{\star}, T^{\star}\right)$.

## §2. Complete Riemannian Manifolds

Let $(M, g)$ be a riemannian manifold of dimension $m$, with metric

$$
g(x)=\sum g_{j k}(x) d x_{j} \otimes d x_{k}, \quad 1 \leqslant j, k \leqslant m
$$

The length of a path $\gamma:[a, b] \longrightarrow M$ is by definition

$$
\ell(\gamma)=\int_{a}^{b}\left|\gamma^{\prime}(t)\right|_{g} d t=\int_{a}^{b}\left(\sum_{j, k} g_{j k}(\gamma(t)) \gamma_{j}^{\prime}(t) \gamma_{k}^{\prime}(t)\right)^{1 / 2} d t
$$

The geodesic distance of two points $x, y \in M$ is

$$
\delta(x, y)=\inf _{\gamma} \ell(\gamma) \quad \text { over paths } \gamma \text { with } \gamma(a)=x, \quad \gamma(b)=y
$$

if $x, y$ are in the same connected component of $M, \delta(x, y)=+\infty$ otherwise. It is easy to check that $\delta$ satisfies the usual axioms of distances: for the separation axiom, use the fact that if $y$ is outside some closed coordinate ball $\bar{B}$ of radius $r$ centered at $x$ and if $g \geqslant c|d x|^{2}$ on $\bar{B}$, then $\delta(x, y) \geqslant c^{1 / 2} r$. In addition, $\delta$ satisfies the axiom:

$$
\begin{equation*}
\text { for every } x, y \in M, \quad \inf _{z \in M} \max \{\delta(x, z), \delta(y, z)\}=\frac{1}{2} \delta(x, y) \text {. } \tag{2.1}
\end{equation*}
$$

In fact for every $\varepsilon>0$ there is a path $\gamma$ such that $\gamma(a)=x, \gamma(b)=y, \ell(\gamma)<\delta(x, y)+\varepsilon$ and we can take $z$ to be at mid-distance between $x$ and $y$ along $\gamma$. A metric space $E$ with a distance $\delta$ satisfying the additional axiom (2.1) will be called a geodesic metric space. It is then easy to see by dichotomy that any two points $x, y \in E$ can be joined by a chain of points $x=x_{0}, x_{1}, \ldots, x_{N}=y$ such that $\delta\left(x_{j}, x_{j+1}\right)<\varepsilon$ and $\sum \delta\left(x_{j}, x_{j+1}\right)<\delta(x, y)+\varepsilon$.
(2.2) Lemma (Hopf-Rinow). Let $(E, \delta)$ be a geodesic metric space. Then the following properties are equivalent:
a) $E$ is locally compact and complete;
b) all closed geodesic balls $\bar{B}\left(x_{0}, r\right)$ are compact.

Proof. Since any Cauchy sequence is bounded, it is immediate that b) implies a). We now check that a) $\Longrightarrow \mathrm{b}$ ). Fix $x_{0}$ and define $R$ to be the supremum of all $r>0$ such that $\bar{B}\left(x_{0}, r\right)$ is compact. Since $E$ is locally compact, we have $R>0$. Suppose that $R<+\infty$. Then $\bar{B}\left(x_{0}, r\right)$ is compact for every $r<R$. Let $y_{\nu}$ be a sequence of points in $\bar{B}\left(x_{0}, R\right)$. Fix an integer $p$. As $\delta\left(x_{0}, y_{\nu}\right) \leqslant R$, axiom (2.1) shows that we can find points $z_{\nu} \in M$ such that $\delta\left(x_{0}, z_{\nu}\right) \leqslant\left(1-2^{-p}\right) R$ and $\delta\left(z_{\nu}, y_{\nu}\right) \leqslant 2^{1-p} R$. Since $\bar{B}\left(x_{0},\left(1-2^{-p}\right) R\right)$ is compact, there is a subsequence $\left(z_{\nu(p, q)}\right)_{q \in \mathbb{N}}$ converging to a limit point $w_{p}$ with $\delta\left(z_{\nu(p, q)}, w_{p}\right) \leqslant 2^{-q}$. We proceed by induction on $p$ and take $\nu(p+1, q)$ to be a subsequence of $\nu(p, q)$. Then

$$
\delta\left(y_{\nu(p, q)}, w_{p}\right) \leqslant \delta\left(y_{\nu(p, q)}, z_{\nu(p, q)}\right)+\delta\left(z_{\nu(p, q)}, w_{p}\right) \leqslant 2^{1-p} R+2^{-q} .
$$

Since $\left(y_{\nu(p+1, q)}\right)$ is a subsequence of $\left(y_{\nu(p, q)}\right)$, we infer that $\delta\left(w_{p}, w_{p+1}\right) \leqslant 32^{-p} R$ by letting $q$ tend to $+\infty$. By the completeness hypothesis, the Cauchy sequence ( $w_{p}$ ) converges to a limit point $w \in M$, and the above inequalities show that $\left(y_{\nu(p, p)}\right)$ converges to $w \in \bar{B}\left(x_{0}, R\right)$. Therefore $\bar{B}\left(x_{0}, R\right)$ is compact. Now, each point $y \in \bar{B}\left(x_{0}, R\right)$ can be covered by a compact ball $\bar{B}\left(y, \varepsilon_{y}\right)$, and the compact set $\bar{B}\left(x_{0}, R\right)$ admits a finite covering by concentric balls $B\left(y_{j}, \varepsilon_{y_{j}} / 2\right)$. Set $\varepsilon=\min \varepsilon_{y_{j}}$. Every point $z \in \bar{B}\left(x_{0}, R+\varepsilon / 2\right)$ is at distance $\leqslant \varepsilon / 2$ of some point $y \in \bar{B}\left(x_{0}, R\right)$, hence at distance $\leqslant \varepsilon / 2+\varepsilon_{y_{j}} / 2$ of some point $y_{j}$, in particular $\bar{B}\left(x_{0}, R+\varepsilon / 2\right) \subset \bigcup \bar{B}\left(y_{j}, \varepsilon_{y_{j}}\right)$ is compact. This is a contradiction, so $R=+\infty$.

The following standard definitions and properties will be useful in order to deal with the completeness of the metric.

## (2.3) Definitions.

a) A riemannian manifold $(M, g)$ is said to be complete if $(M, \delta)$ is complete as a metric space.
b) A continuous function $\psi: M \rightarrow \mathbb{R}$ is said to be exhaustive if for every $c \in \mathbb{R}$ the sublevel set $M_{c}=\{x \in M ; \psi(x)<c\}$ is relatively compact in $M$.
c) A sequence $\left(K_{\nu}\right)_{\nu \in \mathbb{N}}$ of compact subsets of $M$ is said to be exhaustive if $M=\bigcup K_{\nu}$ and if $K_{\nu}$ is contained in the interior of $K_{\nu+1}$ for all $\nu$ (so that every compact subset of $M$ is contained in some $K_{\nu}$ ).
(2.4) Lemma. The following properties are equivalent:
a) $(M, g)$ is complete;
b) there exists an exhaustive function $\psi \in \mathscr{C}^{\infty}(M, \mathbb{R})$ such that $|d \psi|_{g} \leqslant 1$;
c) there exists an exhaustive sequence $\left(K_{\nu}\right)_{\nu \in \mathbb{N}}$ of compact subsets of $M$ and functions $\psi_{\nu} \in \mathscr{C}^{\infty}(M, \mathbb{R})$ such that

$$
\begin{aligned}
& \psi_{\nu}=1 \quad \text { in a neighborhood of } K_{\nu}, \quad \text { Supp } \psi_{\nu} \subset K_{\nu+1}^{\circ}, \\
& 0 \leqslant \psi_{\nu} \leqslant 1 \text { and }\left|d \psi_{\nu}\right|_{g} \leqslant 2^{-\nu}
\end{aligned}
$$

Proof. a) $\Longrightarrow \mathrm{b})$. Without loss of generality, we may assume that $M$ is connected. Select a point $x_{0} \in M$ and set $\psi_{0}(x)=\frac{1}{2} \delta\left(x_{0}, x\right)$. Then $\psi_{0}$ is a Lipschitz function with constant $\frac{1}{2}$, thus $\psi_{0}$ is differentiable almost everywhere on $M$ and $\left|d \psi_{0}\right|_{g} \leqslant \frac{1}{2}$. We can find a smoothing $\psi$ of $\psi_{0}$ such that $|d \psi|_{g} \leqslant 1$ and $\left|\psi-\psi_{0}\right| \leqslant 1$. Then $\psi$ is an exhaustion function of $M$.
b) $\Longrightarrow$ c). Choose $\psi$ as in a) and a function $\rho \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$ such that $\rho=1$ on $\left.]-\infty, 1.1\right]$, $\rho=0$ on $\left[1.9,+\infty\left[\right.\right.$ and $0 \leqslant \rho^{\prime} \leqslant 2$ on $[1,2]$. Then

$$
K_{\nu}=\left\{x \in M ; \psi(x) \leqslant 2^{\nu+1}\right\}, \quad \psi_{\nu}(x)=\rho\left(2^{-\nu-1} \psi(x)\right)
$$

satisfy our requirements.
c) $\Longrightarrow \mathrm{b})$. Set $\psi=\sum 2^{\nu}\left(1-\psi_{\nu}\right)$.
b) $\Longrightarrow$ a). The inequality $|d \psi|_{g} \leqslant 1$ implies $|\psi(x)-\psi(y)| \leqslant \delta(x, y)$ for all $x, y \in M$, so all $\delta$-balls must be relatively compact in $M$.

## $\S$ 3. $L^{2}$ Hodge Theory on Complete Riemannian Manifolds

Let $(M, g)$ be a riemannian manifold and let $F_{1}, F_{2}$ be hermitian $\mathscr{C}^{\infty}$ vector bundles over $M$. If $P: \mathscr{C}^{\infty}\left(M, F_{1}\right) \longrightarrow \mathscr{C}^{\infty}\left(M, F_{2}\right)$ is a differential operator with smooth coefficients, then $P$ induces a non bounded operator

$$
\widetilde{P}: L^{2}\left(M, F_{1}\right) \longrightarrow L^{2}\left(M, F_{2}\right)
$$

as follows: if $u \in L^{2}\left(\underset{\sim}{M}, F_{1}\right)$, we compute $\widetilde{P} u$ in the sense of distribution theory and we say that $u \in \operatorname{Dom} \widetilde{P}$ if $\widetilde{P} u \in L^{2}\left(M, F_{2}\right)$. It follows that $\widetilde{P}$ is densely defined, since Dom $P$ contains the set $\mathscr{D}\left(M, F_{1}\right)$ of compactly supported sections of $\mathscr{C}^{\infty}\left(M, F_{1}\right)$, which is dense in $L^{2}\left(M, F_{1}\right)$. Furthermore $\mathrm{Gr} \widetilde{P}$ is closed: if $u_{\nu} \rightarrow u$ in $L^{2}\left(M, F_{1}\right)$ and $\widetilde{P} u_{\nu} \rightarrow v$ in $L^{2}\left(M, F_{2}\right)$ then $\widetilde{P} u_{\nu} \rightarrow \widetilde{P} u$ in the weak topology of distributions, thus we must have $\widetilde{P} u=v$ and $(u, v) \in \operatorname{Gr} \widetilde{P}$. By the general results of $\S 1$, we see that $\widetilde{P}$ has a closed and densely defined Von Neumann adjoint $(\widetilde{P})^{\star}$. We want to stress, however, that $(\widetilde{P})^{\star}$ does not always coincide with the extension $\left(P^{\star}\right)^{\sim}$ of the formal adjoint $P^{\star}$ : $\mathscr{C}^{\infty}\left(M, F_{2}\right) \longrightarrow \mathscr{C}^{\infty}\left(M, F_{1}\right)$, computed in the sense of distribution theory. In fact $u \in$ $\operatorname{Dom}(\widetilde{P})^{\star}$, resp. $u \in \operatorname{Dom}\left(P^{\star}\right)^{\sim}$, if and only if there is an element $v \in L^{2}\left(M, F_{1}\right)$ such that $\langle u, \widetilde{P} f\rangle=\langle v, f\rangle$ for all $f \in \operatorname{Dom} \widetilde{P}$, resp. for all $f \in \mathscr{D}\left(M, F_{1}\right)$. Therefore we always have $\operatorname{Dom}(\widetilde{P})^{\star} \subset \operatorname{Dom}\left(P^{\star}\right)^{\sim}$ and the inclusion may be strict because the integration by parts to perform may involve boundary integrals for $(\widetilde{P})^{\star}$.
(3.1) Example. Consider

$$
P=\frac{d}{d x}: L^{2}(] 0,1[) \longrightarrow L^{2}(] 0,1[)
$$

where the $L^{2}$ space is taken with respect to the Lebesgue measure $d x$. Then $\operatorname{Dom} \widetilde{P}$ consists of all $L^{2}$ functions with $L^{2}$ derivatives on $] 0,1[$. Such functions have a continuous extension to the interval $[0,1]$. An integration by parts shows that

$$
\int_{0}^{1} u \frac{\overline{d f}}{d x} d x=\int_{0}^{1}-\frac{d u}{d x} \bar{f} d x
$$

for all $f \in \mathscr{D}(] 0,1[)$, thus $P^{\star}=-d / d x=-P$. However for $f \in \operatorname{Dom} \widetilde{P}$ the integration by parts involves the extra term $u(1) \bar{f}(1)-u(0) \bar{f}(0)$ in the right hand side, which is thus continuous in $f$ with respect to the $L^{2}$ topology if and only if $d u / d x \in L^{2}$ and $u(0)=u(1)=0$. Therefore $\operatorname{Dom}(\widetilde{P})^{\star}$ consists of all $u \in \operatorname{Dom}\left(P^{\star}\right)^{\sim}=\operatorname{Dom} \widetilde{P}$ satisfying the additional boundary condition $u(0)=u(1)=0$.

Let $E \rightarrow M$ be a differentiable hermitian bundle. In what follows, we still denote by $D, \delta, \Delta$ the differential operators of $\S$ VI-2 extended in the sense of distribution theory (as explained above). These operators are thus closed and densely defined operators on $L_{\bullet}^{2}(M, E)=\bigoplus_{p} L_{p}^{2}(M, E)$. We also introduce the space $\mathscr{D}_{p}(M, E)$ of compactly supported forms in $\mathscr{C}_{p}^{\infty}(M, E)$. The theory relies heavily on the following important result.
(3.2) Theorem. Assume that $(M, g)$ is complete. Then
a) $\mathscr{D}_{\bullet}(M, E)$ is dense in $\operatorname{Dom} D, \operatorname{Dom} \delta$ and $\operatorname{Dom} D \cap \operatorname{Dom} \delta$ respectively for the graph norms

$$
u \mapsto\|u\|+\|D u\|, \quad u \mapsto\|u\|+\|\delta u\|, \quad u \mapsto\|u\|+\|D u\|+\|\delta u\| .
$$

b) $D^{\star}=\delta, \delta^{\star}=D$ as adjoint operators in Von Neumann's sense.
c) One has $\langle u, \Delta u\rangle=\|D u\|^{2}+\|\delta u\|^{2}$ for every $u \in \operatorname{Dom} \Delta$. In particular
$\operatorname{Dom} \Delta \subset \operatorname{Dom} D \cap \operatorname{Dom} \delta, \quad \operatorname{Ker} \Delta=\operatorname{Ker} D \cap \operatorname{Ker} \delta$,
and $\Delta$ is self-adjoint.
d) If $D^{2}=0$, there are orthogonal decompositions

$$
\begin{aligned}
L_{\bullet}^{2}(M, E) & =\mathscr{H}^{\bullet}(M, E) \oplus \overline{\operatorname{Im} D} \oplus \overline{\operatorname{Im} \delta}, \\
\text { Ker } D & =\mathscr{H}^{\bullet}(M, E) \oplus \overline{\operatorname{Im} D},
\end{aligned}
$$

where $\mathscr{H}^{\bullet}(M, E)=\left\{u \in L_{\bullet}^{2}(M, E) ; \Delta u=0\right\} \subset \mathscr{C}_{\bullet}^{\infty}(M, E)$ is the space of $L^{2}$ harmonic forms.

Proof. a) We show that every element $u \in \operatorname{Dom} D$ can be approximated in the graph norm of $D$ by smooth and compactly supported forms. By hypothesis, $u$ and $D u$ belong to $L_{\bullet}^{2}(M, E)$. Let $\left(\psi_{\nu}\right)$ be a sequence of functions as in Lemma 2.4 c$)$. Then $\psi_{\nu} u \rightarrow u$ in $L_{\bullet}^{2}(M, E)$ and $D\left(\psi_{\nu} u\right)=\psi_{\nu} D u+d \psi_{\nu} \wedge u$ where

$$
\left|d \psi_{\nu} \wedge u\right| \leqslant\left|d \psi_{\nu}\right||u| \leqslant 2^{-\nu}|u| .
$$

Therefore $d \psi_{\nu} \wedge u \rightarrow 0$ and $D\left(\psi_{\nu} u\right) \rightarrow D u$. After replacing $u$ by $\psi_{\nu} u$, we may assume that $u$ has compact support, and by using a finite partition of unity on a neighborhood of Supp $u$ we may also assume that Supp $u$ is contained in a coordinate chart of $M$ on which $E$ is trivial. Let $A$ be the connection form of $D$ on this chart and $\left(\rho_{\varepsilon}\right)$ a family of smoothing kernels. Then $u \star \rho_{\varepsilon} \in \mathscr{D} \bullet(M, E)$ converges to $u$ in $L^{2}(M, E)$ and

$$
D\left(u \star \rho_{\varepsilon}\right)-(D u) \star \rho_{\varepsilon}=A \wedge\left(u \star \rho_{\varepsilon}\right)-(A \wedge u) \star \rho_{\varepsilon}
$$

because $d$ commutes with convolution (as any differential operator with constant coefficients). Moreover $(D u) \star \rho_{\varepsilon}$ converges to $D u$ in $L^{2}(M, E)$ and $A \wedge\left(u \star \rho_{\varepsilon}\right),(A \wedge u) \star \rho_{\varepsilon}$ both converge to $A \wedge u$ since $A \wedge \bullet$ acts continuously on $L^{2}$. Thus $D\left(u \star \rho_{\varepsilon}\right)$ converges to $D u$ and the density of $\mathscr{D}_{\bullet}(M, E)$ in $\operatorname{Dom} D$ follows. The proof for $\operatorname{Dom} \delta$ and $\operatorname{Dom} D \cap \operatorname{Dom} \delta$ is similar, except that the principal part of $\delta$ no longer has constant coefficients in general. The convolution technique requires in this case the following lemma due to K.O. Friedrichs (see e.g. [Hörmander 1963]).
(3.3) Lemma. Let $P f=\sum a_{k} \partial f / \partial x_{k}+b f$ be a differential operator of order 1 on an open set $\Omega \subset \mathbb{R}^{n}$, with coefficients $a_{k} \in C^{1}(\Omega), b \in C^{0}(\Omega)$. Then for any $v \in L^{2}\left(\mathbb{R}^{n}\right)$ with compact support in $\Omega$ we have

$$
\lim _{\varepsilon \rightarrow 0}\left\|P\left(v \star \rho_{\varepsilon}\right)-(P v) \star \rho_{\varepsilon}\right\|_{L^{2}}=0
$$

Proof. It is enough to consider the case when $P=a \partial / \partial x_{k}$. As the result is obvious if $v \in C^{1}$, we only have to show that

$$
\left\|P\left(v \star \rho_{\varepsilon}\right)-(P v) \star \rho_{\varepsilon}\right\|_{L^{2}} \leqslant C\|v\|_{L^{2}}
$$

and to use a density argument. A computation of $w_{\varepsilon}=P\left(v \star \rho_{\varepsilon}\right)-(P v) \star \rho_{\varepsilon}$ by means of an integration by parts gives

$$
\begin{gathered}
w_{\varepsilon}(x)=\int_{\mathbb{R}^{n}}\left(a(x) \frac{\partial v}{\partial x_{k}}(x-\varepsilon y) \rho(y)-a(x-\varepsilon y) \frac{\partial v}{\partial x_{k}}(x-\varepsilon y) \rho(y)\right) d y \\
=\int_{\mathbb{R}^{n}}\left((a(x)-a(x-\varepsilon y)) v(x-\varepsilon y) \frac{1}{\varepsilon} \partial_{k} \rho(y)\right. \\
\left.\quad+\partial_{k} a(x-\varepsilon y) v(x-\varepsilon y) \rho(y)\right) d y
\end{gathered}
$$

If $C$ is a bound for $|d a|$ in a neighborhood of Supp $v$, we get

$$
\left|w_{\varepsilon}(x)\right| \leqslant C \int_{\mathbb{R}^{n}}|v(x-\varepsilon y)|\left(|y|\left|\partial_{k} \rho(y)\right|+|\rho(y)|\right) d y
$$

so Minkowski's inequality $\|f \star g\|_{L^{p}} \leqslant\|f\|_{L^{1}}\|g\|_{L^{p}}$ gives

$$
\left\|w_{\varepsilon}\right\|_{L^{2}} \leqslant C\left(\int_{\mathbb{R}^{n}}\left(|y|\left|\partial_{k} \rho(y)\right|+|\rho(y)|\right) d y\right)\|v\|_{L^{2}}
$$

Proof (end).. b) is equivalent to the fact that

$$
\langle\langle D u, v\rangle\rangle=\langle\langle u, \delta v\rangle\rangle, \quad \forall u \in \operatorname{Dom} D, \quad \forall v \in \operatorname{Dom} \delta
$$

By a), we can find $u_{\nu}, v_{\nu} \in \mathscr{D}_{\bullet}(M, E)$ such that

$$
u_{\nu} \rightarrow u, \quad v_{\nu} \rightarrow v, \quad D u_{\nu} \rightarrow D u \quad \text { and } \quad \delta v_{\nu} \rightarrow \delta v \quad \text { in } \quad L_{\bullet}^{2}(M, E)
$$

and the required equality is the limit of the equalities $\left\langle\left\langle D u_{\nu}, v_{\nu}\right\rangle\right\rangle=\left\langle\left\langle u_{\nu}, \delta v_{\nu}\right\rangle\right\rangle$.
c) Let $u \in \operatorname{Dom} \Delta$. As $\Delta$ is an elliptic operator of order $2, u$ must be in $W_{\bullet}^{2}(M, E, \operatorname{loc})$ by Gårding's inequality. In particular $D u, \delta u \in L^{2}(M, E, l o c)$ and we can perform all integrations by parts that we want if the forms are multiplied by compactly supported functions $\psi_{\nu}$. Let us compute

$$
\begin{aligned}
& \left\|\psi_{\nu} D u\right\|^{2}+\left\|\psi_{\nu} \delta u\right\|^{2}= \\
& =\left\langle\left\langle\psi_{\nu}^{2} D u, D u\right\rangle\right\rangle+\left\langle\left\langle u, D\left(\psi_{\nu}^{2} \delta u\right)\right\rangle\right\rangle \\
& =\left\langle\left\langle D\left(\psi_{\nu}^{2} u\right), D u\right\rangle\right\rangle+\left\langle\left\langle u, \psi_{\nu}^{2} D \delta u\right\rangle-2\left\langle\left\langle\psi_{\nu} d \psi_{\nu} \wedge u, D u\right\rangle\right\rangle+2\left\langle\left\langle u, \psi_{\nu} d \psi_{\nu} \wedge \delta u\right\rangle\right\rangle\right. \\
& =\left\langle\left\langle\psi_{\nu}^{2} u, \Delta u\right\rangle\right\rangle-2\left\langle\left\langle d \psi_{\nu} \wedge u, \psi_{\nu} D u\right\rangle\right\rangle+2\left\langle\left\langle u, d \psi_{\nu} \wedge\left(\psi_{\nu} \delta u\right)\right\rangle\right\rangle \\
& \leqslant\left\langle\left\langle\psi_{\nu}^{2} u, \Delta u\right\rangle\right\rangle+2^{-\nu}\left(2\left\|\psi_{\nu} D u\right\|\|u\|+2\left\|\psi_{\nu} \delta u\right\|\|u\|\right) \\
& \leqslant\left\langle\left\langle\psi_{\nu}^{2} u, \Delta u\right\rangle\right\rangle+2^{-\nu}\left(\left\|\psi_{\nu} D u\right\|^{2}+\left\|\psi_{\nu} \delta u\right\|^{2}+2\|u\|^{2}\right) .
\end{aligned}
$$

We get therefore

$$
\left\|\psi_{\nu} D u\right\|^{2}+\left\|\psi_{\nu} \delta u\right\|^{2} \leqslant \frac{1}{1-2^{-\nu}}\left(\left\langle\left\langle\psi_{\nu}^{2} u, \Delta u\right\rangle+2^{1-\nu}\|u\|^{2}\right) .\right.
$$

By letting $\nu$ tend to $+\infty$, we obtain $\|D u\|^{2}+\|\delta u\|^{2} \leqslant\langle\langle u, \Delta u\rangle\rangle$, in particular $D u, \delta u$ are in $L_{\bullet}^{2}(M, E)$. This implies

$$
\langle\langle u, \Delta v\rangle\rangle=\langle\langle D u, D v\rangle\rangle+\langle\langle\delta u, \delta v\rangle\rangle, \quad \forall u, v \in \operatorname{Dom} \Delta,
$$

because the equality holds for $\psi_{\nu} u$ and $v$, and because we have $\psi_{\nu} u \rightarrow u, D\left(\psi_{\nu} u\right) \rightarrow D u$ and $\delta\left(\psi_{\nu} u\right) \rightarrow \delta u$ in $L^{2}$. Therefore $\Delta$ is self-adjoint.
d) is an immediate consequence of b), c) and Th. 1.2.

On a complete hermitian manifold $(X, \omega)$, there are of course similar results for the operators $D^{\prime}, D^{\prime \prime}, \delta^{\prime}, \delta^{\prime \prime}, \Delta^{\prime}, \Delta^{\prime \prime}$ attached to a hermitian vector bundle $E$.

## $\S$ 4. General Estimate for $\boldsymbol{d}^{\prime \prime}$ on Hermitian Manifolds

Let $(X, \omega)$ be a complete hermitian manifold and $E$ a hermitian holomorphic vector bundle of rank $r$ over $X$. Assume that the hermitian operator

$$
\begin{equation*}
A_{E, \omega}=[\mathrm{i} \Theta(E), \Lambda]+T_{\omega} \tag{4.1}
\end{equation*}
$$

is semi-positive on $\Lambda^{p, q} T_{X}^{\star} \otimes E$. Then for every form $u \in \operatorname{Dom} D^{\prime \prime} \cap \operatorname{Dom} \delta^{\prime \prime}$ of bidegree ( $p, q$ ) we have

$$
\begin{equation*}
\left\|D^{\prime \prime} u\right\|^{2}+\left\|\delta^{\prime \prime} u\right\|^{2} \geqslant \int_{X}\left\langle A_{E, \omega} u, u\right\rangle d V \tag{4.2}
\end{equation*}
$$

In fact (4.2) is true for all $u \in \mathscr{D}_{p, q}(X, E)$ in view of the Bochner-Kodaira-Nakano identity VII-2.3, and this result is easily extended to every $u$ in $\operatorname{Dom} D^{\prime \prime} \cap \operatorname{Dom} \delta^{\prime \prime}$ by density of $\mathscr{D}_{p, q}(X, E)(T h .3 .2$ a) $)$.

Assume now that a form $g \in L_{p, q}^{2}(X, E)$ is given such that

$$
\begin{equation*}
D^{\prime \prime} g=0 \tag{4.3}
\end{equation*}
$$

and that for almost every $x \in X$ there exists $\alpha \in[0,+\infty[$ such that

$$
|\langle g(x), u\rangle|^{2} \leqslant \alpha\left\langle A_{E, \omega} u, u\right\rangle
$$

for every $u \in\left(\Lambda^{p, q} T_{X}^{\star} \otimes E\right)_{x}$. If the operator $A_{E, \omega}$ is invertible, the minimal such number $\alpha$ is $\left|A_{E, \omega}^{-1 / 2} g(x)\right|^{2}=\left\langle A_{E, \omega}^{-1} g(x), g(x)\right\rangle$, so we shall always denote it in this way even when $A_{E, \omega}$ is no longer invertible. Assume furthermore that

$$
\begin{equation*}
\int_{X}\left\langle A_{E, \omega}^{-1} g, g\right\rangle d V<+\infty \tag{4.4}
\end{equation*}
$$

The basic result of $L^{2}$ theory can be stated as follows.
(4.5) Theorem. If $(X, \omega)$ is complete and $A_{E, \omega} \geqslant 0$ in bidegree $(p, q)$, then for any $g \in L_{p, q}^{2}(X, E)$ satisfying (4.4) such that $D^{\prime \prime} g=0$ there exists $f \in L_{p, q-1}^{2}(X, E)$ such that $D^{\prime \prime} f=g$ and

$$
\|f\|^{2} \leqslant \int_{X}\left\langle A_{E, \omega}^{-1} g, g\right\rangle d V
$$

Proof. For every $u \in \operatorname{Dom} D^{\prime \prime} \cap \operatorname{Dom} \delta^{\prime \prime}$ we have

$$
\begin{aligned}
|\langle\langle u, g\rangle\rangle|^{2}=\left|\int_{X}\langle u, g\rangle d V\right|^{2} & \leqslant\left(\int_{X}\left\langle A_{E, \omega} u, u\right\rangle^{1 / 2}\left\langle A_{E, \omega}^{-1} g, g\right\rangle^{1 / 2} d V\right)^{2} \\
& \leqslant \int_{X}\left\langle A_{E, \omega}^{-1} g, g\right\rangle d V \cdot \int_{X}\left\langle A_{E, \omega} u, u\right\rangle d V
\end{aligned}
$$

by means of the Cauchy-Schwarz inequality. The a priori estimate (4.2) implies

$$
|\langle\langle u, g\rangle\rangle|^{2} \leqslant C\left(\left\|D^{\prime \prime} u\right\|^{2}+\left\|\delta^{\prime \prime} u\right\|^{2}\right), \quad \forall u \in \operatorname{Dom} D^{\prime \prime} \cap \operatorname{Dom} \delta^{\prime \prime}
$$

where $C$ is the integral (4.4). Now we just have to repeat the proof of the existence part of Th. 1.2. For any $u \in \operatorname{Dom} \delta^{\prime \prime}$, let us write

$$
u=u_{1}+u_{2}, \quad u_{1} \in \operatorname{Ker} D^{\prime \prime}, \quad u_{2} \in\left(\operatorname{Ker} D^{\prime \prime}\right)^{\perp}=\overline{\operatorname{Im} \delta^{\prime \prime}}
$$

Then $D^{\prime \prime} u_{1}=0$ and $\delta^{\prime \prime} u_{2}=0$. Since $g \in \operatorname{Ker} D^{\prime \prime}$, we get

$$
\left.|\langle u u, g\rangle\rangle\right|^{2}=\left|\left\langle\left\langle u_{1}, g\right\rangle\right\rangle\right|^{2} \leqslant C\left\|\delta^{\prime \prime} u_{1}\right\|^{2}=C\left\|\delta^{\prime \prime} u\right\|^{2}
$$

The Hahn-Banach theorem shows that the continuous linear form

$$
L_{p, q-1}^{2}(X, E) \ni \delta^{\prime \prime} u \longmapsto\langle\langle u, g\rangle\rangle
$$

can be extended to a linear form $v \longmapsto\langle\langle v, f\rangle\rangle, f \in L_{p, q-1}^{2}(X, E)$, of norm $\|f\| \leqslant C^{1 / 2}$. This means that

$$
\langle\langle u, g\rangle\rangle=\left\langle\left\langle\delta^{\prime \prime} u, f\right\rangle\right\rangle, \quad \forall u \in \operatorname{Dom} \delta^{\prime \prime}
$$

i.e. that $D^{\prime \prime} f=g$. The theorem is proved.
(4.6) Remark. One can always find a solution $f \in\left(\operatorname{Ker} D^{\prime \prime}\right)^{\perp}$ : otherwise replace f by its orthogonal projection on $\left(\operatorname{Ker} D^{\prime \prime}\right)^{\perp}$. This solution is clearly unique and is precisely the solution of minimal $L^{2}$ norm of the equation $D^{\prime \prime} f=g$. We have $f \in \overline{\overline{\operatorname{Im}} \delta^{\prime \prime}}$, thus $f$ satifies the additional equation

$$
\begin{equation*}
\delta^{\prime \prime} f=0 . \tag{4.7}
\end{equation*}
$$

Consequently $\Delta^{\prime \prime} f=\delta^{\prime \prime} D^{\prime \prime} f=\delta^{\prime \prime} g$. If $g \in \mathscr{C}_{p, q}^{\infty}(X, E)$, the ellipticity of $\Delta^{\prime \prime}$ shows that $f \in \mathscr{C}_{p, q-1}^{\infty}(X, E)$.
(4.8) Remark. If $A_{E, \omega}$ is positive definite, let $\lambda(x)>0$ be the smallest eigenvalue of this operator at $x \in X$. Then $\lambda$ is continuous on $X$ and we have

$$
\int_{X}\left\langle A_{E, \omega}^{-1} g, g\right\rangle d V \leqslant \int_{X} \lambda(x)^{-1}|g(x)|^{2} d V .
$$

The above situation occurs for example if $\omega$ is complete Kähler, $E>_{m} 0$ and $p=n$, $q \geqslant 1, m \geqslant \min \{n-q+1, r\}$ (apply Lemma VII-7.2).

## § 5. Estimates on Weakly Pseudoconvex Manifolds

We first introduce a large class of complex manifolds on which the $L^{2}$ estimates will be easily tractable.
(5.1) Definition. A complex manifold $X$ is said to be weakly pseudoconvex if there exists an exhaustion function $\psi \in \mathscr{C}^{\infty}(X, \mathbb{R})$ such that $\mathrm{id} d^{\prime} d^{\prime \prime} \psi \geqslant 0$ on $X$, i.e. $\psi$ is plurisubharmonic.

For domains $\Omega \subset \mathbb{C}^{n}$, the above weak pseudoconvexity notion is equivalent to pseudoconvexity (cf. Th. I-4.14). Note that every compact manifold is also weakly pseudoconvex (take $\psi \equiv 0$ ). Other examples that will appear later are Stein manifolds, or the total space of a Griffiths semi-negative vector bundle over a compact manifold (cf. Prop. IX?.?).
(5.2) Theorem. Every weakly pseudoconvex Kähler manifold $(X, \omega)$ carries a complete Kähler metric $\widehat{\omega}$.

Proof. Let $\psi \in \mathscr{C}^{\infty}(X, \mathbb{R})$ be an exhaustive plurisubharmonic function on $X$. After addition of a constant to $\psi$, we can assume $\psi \geqslant 0$. Then $\widehat{\omega}=\omega+\mathrm{i} d^{\prime} d^{\prime \prime}\left(\psi^{2}\right)$ is Kähler and

$$
\widehat{\omega}=\omega+2 i \psi d^{\prime} d^{\prime \prime} \psi+2 \mathrm{i} d^{\prime} \psi \wedge d^{\prime \prime} \psi \geqslant \omega+2 \mathrm{i} d^{\prime} \psi \wedge d^{\prime \prime} \psi .
$$

Since $d \psi=d^{\prime} \psi+d^{\prime \prime} \psi$, we get $|d \psi|_{\widehat{\omega}}=\sqrt{2}\left|d^{\prime} \psi\right|_{\widehat{\omega}} \leqslant 1$ and Lemma 2.4 shows that $\widehat{\omega}$ is complete.

Observe that we could have set more generally $\widehat{\omega}=\omega+\mathrm{i} d^{\prime} d^{\prime \prime}(\chi \circ \psi)$ where $\chi$ is a convex increasing function. Then

$$
\begin{align*}
\widehat{\omega} & =\omega+i\left(\chi^{\prime} \circ \psi\right) d^{\prime} d^{\prime \prime} \psi+i\left(\chi^{\prime \prime} \circ \psi\right) d^{\prime} \psi \wedge d^{\prime \prime} \psi \\
& \geqslant \omega+\mathrm{i} d^{\prime}(\rho \circ \psi) \wedge d^{\prime \prime}(\rho \circ \psi) \tag{5.3}
\end{align*}
$$

where $\rho(t)=\int_{0}^{t} \sqrt{\chi^{\prime \prime}(u)} d u$. We thus have $\left|d^{\prime}(\rho \circ \psi)\right|_{\widehat{\omega}} \leqslant 1$ and $\widehat{\omega}$ will be complete as soon as $\lim _{t \rightarrow+\infty} \rho(t)=+\infty$, i.e.

$$
\begin{equation*}
\int_{0}^{+\infty} \sqrt{\chi^{\prime \prime}(u)} d u=+\infty \tag{5.4}
\end{equation*}
$$

One can take for example $\chi(t)=t-\log (t)$ for $t \geqslant 1$.
It follows from the above considerations that almost all vanishing theorems for positive vector bundles over compact manifolds are also valid on weakly pseudoconvex manifolds. Let us mention here the analogues of some results proved in Chapter 7.
(5.5) Theorem. For any m-positive vector bundle of rank $r$ over a weakly pseudoconvex manifold $X$, we have $H^{n, q}(X, E)=0$ for all $q \geqslant 1$ and $m \geqslant \min \{n-q+1, r\}$.

Proof. The curvature form $\mathrm{i} \Theta(\operatorname{det} E)$ is a Kähler metric on $X$, hence $X$ possesses a complete Kähler metric $\omega$. Let $\psi \in \mathscr{C}^{\infty}(X, \mathbb{R})$ be an exhaustive plurisubharmonic function. For any convex increasing function $\chi \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$, we denote by $E_{\chi}$ the holomorphic vector bundle $E$ together with the modified metric $|u|_{\chi}^{2}=|u|^{2} \exp (-\chi \circ \psi(x)), u \in E_{x}$. We get

$$
\mathrm{i} \Theta\left(E_{\chi}\right)=\mathrm{i} \Theta(E)+\mathrm{i} d^{\prime} d^{\prime \prime}(\chi \circ \psi) \otimes \operatorname{Id}_{E} \geqslant_{m} \mathrm{i} \Theta(E)
$$

thus $A_{E_{\chi}, \omega} \geqslant A_{E, \omega}>0$ in bidegree $(n, q)$. Let $g$ be a given form of bidegree $(n, q)$ with $L_{\text {loc }}^{2}$ coefficients, such that $D^{\prime \prime} g=0$. The integrals

$$
\int_{X}\left\langle A_{E_{\chi}, \omega}^{-1} g, g\right\rangle_{\chi} d V \leqslant \int_{X}\left\langle A_{E, \omega}^{-1} g, g\right\rangle e^{-\chi \circ \psi} d V, \quad \int_{X}|g|^{2} e^{-\chi \circ \psi} d V
$$

become convergent if $\chi$ grows fast enough. We can thus apply Th. 4.5 to $\left(X, E_{\chi}, \omega\right)$ and find a $(n, q-1)$ form $f$ such that $D^{\prime \prime} f=g$. If $g$ is smooth, Remark 4.6 shows that $f$ can also be chosen smooth.
(5.6) Theorem. If $E$ is a positive line bundle over a weakly pseudoconvex manifold $X$, then $H^{p, q}(X, E)=0$ for $p+q \geqslant n+1$.

Proof. The proof is similar to that of Th. 5.5, except that we use here the Kähler metric

$$
\omega_{\chi}=\mathrm{i} \Theta\left(E_{\chi}\right)=\omega+\mathrm{i} d^{\prime} d^{\prime \prime}(\chi \circ \psi), \quad \omega=\mathrm{i} \Theta(E)
$$

which depends on $\chi$. $\operatorname{By}(5.4) \omega_{\chi}$ is complete as soon as $\chi$ is a convex increasing function that grows fast enough. Apply now Th. 4.5 to $\left(X, E_{\chi}, \omega_{\chi}\right)$ and observe that $A_{E_{\chi}, \omega_{\chi}}=$ $\left[\mathrm{i} \Theta\left(E_{\chi}\right), \Lambda_{\chi}\right]=(p+q-n)$ Id in bidegree $(p, q)$ in virtue of Cor. VI-8.4 It remains to show that for every form $g \in \mathscr{C}_{p, q}^{\infty}(X, E)$ there exists a choice of $\chi$ such that $g \in$ $L_{p, q}^{2}\left(X, E_{\chi}, \omega_{\chi}\right)$. By (5.3) the norm of a scalar form with respect to $\omega_{\chi}$ is less than its norm with respect to $\omega$, hence $|g|_{\chi}^{2} \leqslant|g|^{2} \exp (-\chi \circ \psi)$. On the other hand

$$
d V_{\chi} \leqslant C\left(1+\chi^{\prime} \circ \psi+\chi^{\prime \prime} \circ \psi\right)^{n} d V
$$

where $C$ is a positive continuous function on $X$. The following lemma implies that we can always choose $\chi$ in order that the integral of $|g|_{\chi}^{2} d V_{\chi}$ converges on $X$.
(5.7) Lemma. For any positive function $\lambda \in \mathscr{C}^{\infty}([0,+\infty[, \mathbb{R})$, there exists a smooth convex function $\chi \in \mathscr{C}^{\infty}\left(\left[0,+\infty[, \mathbb{R})\right.\right.$ such that $\chi, \chi^{\prime}, \chi^{\prime \prime} \geqslant \lambda$ and $\left(1+\chi^{\prime}+\chi^{\prime \prime}\right)^{n} e^{-\chi} \leqslant 1 / \lambda$.

Proof. We shall construct $\chi$ such that $\chi^{\prime \prime} \geqslant \chi^{\prime} \geqslant \chi \geqslant \lambda$ and $\chi^{\prime \prime} / \chi^{2} \leqslant C$ for some constant $C$. Then $\chi$ satisties the conclusion of the lemma after addition of a constant. Without loss of generality, we may assume that $\lambda$ is increasing and $\lambda \geqslant 1$. We define $\chi$ as a power series

$$
\chi(t)=\sum_{k=0}^{+\infty} a_{0} a_{1} \ldots a_{k} t^{k}
$$

where $a_{k}>0$ is a decreasing sequence converging to 0 very slowly. Then $\chi$ is real analytic on $\mathbb{R}$ and the inequalities $\chi^{\prime \prime} \geqslant \chi^{\prime} \geqslant \chi$ are realized if we choose $a_{k} \geqslant 1 / k, k \geqslant 1$. Select a strictly increasing sequence of integers $\left(N_{p}\right)_{p \geqslant 1}$ so large that $\frac{1}{p} \lambda(p+1)^{1 / N_{p}} \in$ $[1 / p, 1 /(p-1)[$. We set

$$
\begin{aligned}
& a_{0}=\ldots=a_{N_{1}-1}=e \lambda(2), \\
& a_{k}=\frac{1}{p} \lambda(p+1)^{1 / N_{p}} e^{1 / \sqrt{k}}, \quad N_{p} \leqslant k<N_{p+1} .
\end{aligned}
$$

Then $\left(a_{k}\right)$ is decreasing. For $t \in[0,1]$ we have $\chi(t) \geqslant a_{0} \geqslant \lambda(t)$ and for $t \in[1,+\infty[$ the choice $k=N_{p}$ where $p=[t]$ is the integer part of $t$ gives

$$
\chi(t) \geqslant \chi(p) \geqslant\left(a_{0} a_{1} \ldots a_{k}\right) p^{k} \geqslant\left(a_{k} p\right)^{k} \geqslant \lambda(p+1) \geqslant \lambda(t) .
$$

Furthermore, we have

$$
\begin{gathered}
\chi(t)^{2} \geqslant \sum_{k \geqslant 0}\left(a_{0} a_{1} \ldots a_{k}\right)^{2} t^{2 k}, \\
\chi^{\prime \prime}(t)=\sum_{k \geqslant 0}(k+1)(k+2) a_{0} a_{1} \ldots a_{k+2} t^{k},
\end{gathered}
$$

thus we will get $\chi^{\prime \prime}(t) \leqslant C \chi(t)^{2}$ if we can prove that

$$
m^{2} a_{0} a_{1} \ldots a_{2 m} \leqslant C^{\prime}\left(a_{0} a_{1} \ldots a_{m}\right)^{2}, \quad m \geqslant 0
$$

However, as $\frac{1}{p} \lambda(p+1)^{1 / N_{p}}$ is decreasing, we find

$$
\begin{aligned}
\frac{a_{0} a_{1} \ldots a_{2 m}}{\left(a_{0} a_{1} \ldots a_{m}\right)^{2}} & =\frac{a_{m+1} \ldots a_{2 m}}{a_{0} a_{1} \ldots a_{m}} \\
& \leqslant \exp \left(\frac{1}{\sqrt{m+1}}+\cdots+\frac{1}{\sqrt{2 m}}-\frac{1}{\sqrt{1}}-\cdots-\frac{1}{\sqrt{m}}+O(1)\right) \\
& \leqslant \exp (2 \sqrt{2 m}-4 \sqrt{m}+O(1)) \leqslant C^{\prime} m^{-2}
\end{aligned}
$$

As a last application, we generalize the Girbau vanishing theorem in the case of weakly pseudoconvex manifolds. This result is due to [Abdelkader 1980] and [Ohsawa 1981]. We present here a simplified proof which appeared in [Demailly 1985].
(5.8) Theorem. Let $(X, \omega)$ be a weakly pseudoconvex Kähler manifold. If $E$ is a semipositive line bundle such that $\mathrm{i} \Theta(E)$ has at least $n-s+1$ positive eigenvalues at every point, then

$$
H^{p, q}(X, E)=0 \quad \text { for } \quad p+q \geqslant n+s
$$

Proof. Let $\chi, \rho \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$ be convex increasing functions to be specified later. We use here the hermitian metric

$$
\begin{aligned}
\alpha & =\mathrm{i} \Theta\left(E_{\chi}\right)+\exp (-\rho \circ \psi) \omega \\
& =\mathrm{i} \Theta(E)+\mathrm{i} d^{\prime} d^{\prime \prime}(\chi \circ \psi)+\exp (-\rho \circ \psi) \omega .
\end{aligned}
$$

Although $\omega$ is Kähler, the metric $\alpha$ is not so. Denote by $\gamma_{j}^{\chi, \omega}$ (resp. $\gamma_{j}^{\chi, \alpha}$ ), $1 \leqslant j \leqslant n$, the eigenvalues of $\mathrm{i} \Theta\left(E_{\chi}\right)$ with respect to $\omega$ (resp. $\alpha$ ), rearranged in increasing order. The minimax principle implies $\gamma_{j}^{\chi, \omega} \geqslant \gamma_{j}^{0, \omega}$, and the hypothesis yields $0<\gamma_{s}^{0, \omega} \leqslant \gamma_{s+1}^{0, \omega} \leqslant$ $\ldots \leqslant \gamma_{n}^{0, \omega}$ on $X$. By means of a diagonalization of $\mathrm{i} \Theta\left(E_{\chi}\right)$ with respect to $\omega$, we find

$$
1 \geqslant \gamma_{j}^{\chi, \alpha}=\frac{\gamma_{j}^{\chi, \omega}}{\gamma_{j}^{\chi, \omega}+\exp (-\rho \circ \psi)} \geqslant \frac{\gamma_{j}^{0, \omega}}{\gamma_{j}^{0, \omega}+\exp (-\rho \circ \psi)} .
$$

Let $\varepsilon>0$ be small. Select $\rho$ such that $\exp (-\rho \circ \psi(x)) \leqslant \varepsilon \gamma_{s}^{0, \omega}(x)$ at every point. Then for $j \geqslant s$ we get

$$
\gamma_{j}^{\chi, \alpha} \geqslant \frac{\gamma_{j}^{0, \omega}}{\gamma_{j}^{0, \omega}+\varepsilon \gamma_{j}^{0, \omega}}=\frac{1}{1+\varepsilon} \geqslant 1-\varepsilon
$$

and Th. VI-8.3 implies

$$
\begin{aligned}
\left\langle\left[\mathrm{i} \Theta\left(E_{\chi}\right), \Lambda_{\alpha}\right] u, u\right\rangle_{\alpha} & \geqslant\left(\gamma_{1}^{\chi, \alpha}+\cdots+\gamma_{p}^{\chi, \alpha}-\gamma_{q+1}^{\chi, \alpha}-\ldots-\gamma_{n}^{\chi, \alpha}\right)|u|^{2} \\
& \geqslant((p-s+1)(1-\varepsilon)-(n-q))|u|^{2} \\
& \geqslant(1-(p-s+1) \varepsilon)|u|^{2} .
\end{aligned}
$$

It remains however to control the torsion term $T_{\alpha}$. As $\omega$ is Kähler, trivial computations yield

$$
\begin{aligned}
d^{\prime} \alpha & =-\rho^{\prime} \circ \psi \exp (-\rho \circ \psi) d^{\prime} \psi \wedge \omega \\
d^{\prime} d^{\prime \prime} \alpha & =\exp (-\rho \circ \psi)\left[\left(\left(\rho^{\prime} \circ \psi\right)^{2}-\rho^{\prime \prime} \circ \psi\right) d^{\prime} \psi \wedge d^{\prime \prime} \psi-\rho^{\prime} \circ \psi d^{\prime} d^{\prime \prime} \psi\right] \wedge \omega
\end{aligned}
$$

Since

$$
\alpha \geqslant i\left(\chi^{\prime} \circ \psi d^{\prime} d^{\prime \prime} \psi+\chi^{\prime \prime} \circ \psi d^{\prime} \psi \wedge d^{\prime \prime} \psi\right)+\exp (-\rho \circ \psi) \omega,
$$

we get the upper bounds

$$
\begin{aligned}
\left|d^{\prime} \alpha\right|_{\alpha} & \leqslant \rho^{\prime} \circ \psi\left|d^{\prime} \psi\right|_{\alpha}|\exp (-\rho \circ \psi) \omega|_{\alpha} \leqslant \rho^{\prime} \circ \psi\left(\chi^{\prime \prime} \circ \psi\right)^{-\frac{1}{2}} \\
\left|d^{\prime} d^{\prime \prime} \alpha\right|_{\alpha} & \leqslant \frac{\left(\rho^{\prime} \circ \psi\right)^{2}+\rho^{\prime \prime} \circ \psi}{\chi^{\prime \prime} \circ \psi}+\frac{\rho^{\prime} \circ \psi}{\chi^{\prime} \circ \psi} .
\end{aligned}
$$

It is then clear that we can choose $\chi$ growing sufficiently fast in order that $\left|T_{\alpha}\right|_{\alpha} \leqslant \varepsilon$. If $\varepsilon$ is chosen sufficiently small, we get $A_{E_{\chi}, \alpha} \geqslant \frac{1}{2} \mathrm{Id}$, and the conclusion is obtained in the same way as for Th. 5.6.

## § 6. Hörmander's Estimates for non Complete Kähler Metrics

Our aim here is to derive also estimates for a non complete Kähler metric, for example the standard metric of $\mathbb{C}^{n}$ on a bounded domain $\Omega \subset \subset \mathbb{C}^{n}$. A result of this type can be obtained in the situation described at the end of Remark 4.8. The underlying idea is due to [Hörmander 1966], although we do not apply his so called "three weights" technique, but use instead an approximation of the given metric $\omega$ by complete Kähler metrics.
(6.1) Theorem. Let $(X, \widehat{\omega})$ be a complete Kähler manifold, $\omega$ another Kähler metric, possibly non complete, and $E \longrightarrow X$ a m-semi-positive vector bundle. Let $g \in L_{n, q}^{2}(X, E)$ be such that $D^{\prime \prime} g=0$ and

$$
\int_{X}\left\langle A_{q}^{-1} g, g\right\rangle d V<+\infty
$$

with respect to $\omega$, where $A_{q}$ stands for the operator $\mathrm{i} \Theta(E) \wedge \Lambda$ in bidegree $(n, q)$ and $q \geqslant 1$, $m \geqslant \min \{n-q+1, r\}$. Then there exists $f \in L_{n, q-1}^{2}(X, E)$ such that $D^{\prime \prime} f=g$ and

$$
\|f\|^{2} \leqslant \int_{X}\left\langle A_{q}^{-1} g, g\right\rangle d V .
$$

Proof. For every $\varepsilon>0$, the Kähler metric

$$
\omega_{\varepsilon}=\omega+\varepsilon \widehat{\omega}
$$

is complete. The idea of the proof is to apply the $L^{2}$ estimates to $\omega_{\varepsilon}$ and to let $\varepsilon$ tend to zero. Let us put an index $\varepsilon$ to all objects depending on $\omega_{\varepsilon}$. It follows from Lemma 6.3 below that

$$
\begin{equation*}
|u|_{\varepsilon}^{2} d V_{\varepsilon} \leqslant|u|^{2} d V, \quad\left\langle A_{q, \varepsilon}^{-1} u, u\right\rangle_{\varepsilon} d V_{\varepsilon} \leqslant\left\langle A_{q}^{-1} u, u\right\rangle d V \tag{6.2}
\end{equation*}
$$

for every $u \in \Lambda^{n, q} T_{X}^{\star} \otimes E$. If these estimates are taken for granted, Th. 4.5 applied to $\omega_{\varepsilon}$ yields a section $f_{\varepsilon} \in L_{n, q-1}^{2}(X, E)$ such that $D^{\prime \prime} f_{\varepsilon}=g$ and

$$
\int_{X}\left|f_{\varepsilon}\right|_{\varepsilon}^{2} d V_{\varepsilon} \leqslant \int_{X}\left\langle A_{q, \varepsilon}^{-1} g, g\right\rangle_{\varepsilon} d V_{\varepsilon} \leqslant \int_{X}\left\langle A_{q}^{-1} g, g\right\rangle d V .
$$

This implies that the family $\left(f_{\varepsilon}\right)$ is bounded in $L^{2}$ norm on every compact subset of $X$. We can thus find a weakly convergent subsequence $\left(f_{\varepsilon_{\nu}}\right)$ in $L_{\text {loc }}^{2}$. The weak limit $f$ is the solution we are looking for.
(6.3) Lemma. Let $\omega$, $\gamma$ be hermitian metrics on $X$ such that $\gamma \geqslant \omega$. For every $u \in \Lambda^{n, q} T_{X}^{\star} \otimes E, q \geqslant 1$, we have

$$
|u|_{\gamma}^{2} d V_{\gamma} \leqslant|u|^{2} d V, \quad\left\langle A_{q, \gamma}^{-1} u, u\right\rangle_{\gamma} d V_{\gamma} \leqslant\left\langle A_{q}^{-1} u, u\right\rangle d V
$$

where an index $\gamma$ means that the corresponding term is computed in terms of $\gamma$ instead of $\omega$.

Proof. Let $x_{0} \in X$ be a given point and $\left(z_{1}, \ldots, z_{n}\right)$ coordinates such that

$$
\omega=\mathrm{i} \sum_{1 \leqslant j \leqslant n} d z_{j} \wedge d \bar{z}_{j}, \quad \gamma=\mathrm{i} \sum_{1 \leqslant j \leqslant n} \gamma_{j} d z_{j} \wedge d \bar{z}_{j} \quad \text { at } \quad x_{0},
$$

where $\gamma_{1} \leqslant \ldots \leqslant \gamma_{n}$ are the eigenvalues of $\gamma$ with respect to $\omega$ (thus $\gamma_{j} \geqslant 1$ ). We have $\left|d z_{j}\right|_{\gamma}^{2}=\gamma_{j}^{-1}$ and $\left|d z_{K}\right|_{\gamma}^{2}=\gamma_{K}^{-1}$ for any multi-index $K$, with the notation $\gamma_{K}=\prod_{j \in K} \gamma_{j}$. For every $u=\sum u_{K, \lambda} d z_{1} \wedge \ldots \wedge d z_{n} \wedge d \bar{z}_{K} \otimes e_{\lambda},|K|=q, 1 \leqslant \lambda \leqslant r$, the computations of § VII-7 yield

$$
\begin{aligned}
|u|_{\gamma}^{2} & =\sum_{K, \lambda}\left(\gamma_{1} \ldots \gamma_{n}\right)^{-1} \gamma_{K}^{-1}\left|u_{K, \lambda}\right|^{2}, \quad d V_{\gamma}=\gamma_{1} \ldots \gamma_{n} d V, \\
|u|_{\gamma}^{2} d V_{\gamma} & =\sum_{K, \lambda} \gamma_{K}^{-1}\left|u_{K, \lambda}\right|^{2} d V \leqslant|u|^{2} d V, \\
\Lambda_{\gamma} u & =\sum_{|I|=q-1} \sum_{j, \lambda} \mathrm{i}(-1)^{n+j-1} \gamma_{j}^{-1} u_{j I, \lambda}\left(\widehat{d z_{j}}\right) \wedge d \bar{z}_{I} \otimes e_{\lambda},
\end{aligned}
$$

where ( $\widehat{d z_{j}}$ ) means $d z_{1} \wedge \ldots \widehat{d z_{j}} \ldots \wedge d z_{n}$,

$$
\begin{aligned}
A_{q, \gamma} u & =\sum_{|I|=q-1} \sum_{j, k, \lambda, \mu} \gamma_{j}^{-1} c_{j k \lambda \mu} u_{j I, \lambda} d z_{1} \wedge \ldots \wedge d z_{n} \wedge d \bar{z}_{k I} \otimes e_{\mu}, \\
\left\langle A_{q, \gamma} u, u\right\rangle_{\gamma} & =\left(\gamma_{1} \ldots \gamma_{n}\right)^{-1} \sum_{|I|=q-1} \gamma_{I}^{-1} \sum_{j, k, \lambda, \mu} \gamma_{j}^{-1} \gamma_{k}^{-1} c_{j k \lambda \mu} u_{j I, \lambda} \bar{u}_{k I, \mu} \\
& \geqslant\left(\gamma_{1} \ldots \gamma_{n}\right)^{-1} \sum_{|I|=q-1} \gamma_{I}^{-2} \sum_{j, k, \lambda, \mu} \gamma_{j}^{-1} \gamma_{k}^{-1} c_{j k \lambda \mu} u_{j I, \lambda} \bar{u}_{k I, \mu} \\
& =\gamma_{1} \ldots \gamma_{n}\left\langle A_{q} S_{\gamma} u, S_{\gamma} u\right\rangle
\end{aligned}
$$

where $S_{\gamma}$ is the operator defined by

$$
S_{\gamma} u=\sum_{K}\left(\gamma_{1} \ldots \gamma_{n} \gamma_{K}\right)^{-1} u_{K, \lambda} d z_{1} \wedge \ldots \wedge d z_{n} \wedge d \bar{z}_{K} \otimes e_{\lambda} .
$$

We get therefore

$$
\begin{aligned}
\left|\langle u, v\rangle_{\gamma}\right|^{2}=\left|\left\langle u, S_{\gamma} v\right\rangle\right|^{2} & \leqslant\left\langle A_{q}^{-1} u, u\right\rangle\left\langle A_{q} S_{\gamma} v, S_{\gamma} v\right\rangle \\
& \leqslant\left(\gamma_{1} \ldots \gamma_{n}\right)^{-1}\left\langle A_{q}^{-1} u, u\right\rangle\left\langle A_{q, \gamma} v, v\right\rangle_{\gamma}
\end{aligned}
$$

and the choice $v=A_{q, \gamma}^{-1} u$ implies

$$
\left\langle A_{q, \gamma}^{-1} u, u\right\rangle_{\gamma} \leqslant\left(\gamma_{1} \ldots \gamma_{n}\right)^{-1}\left\langle A_{q}^{-1} u, u\right\rangle ;
$$

this relation is equivalent to the last one in the lemma.
An important special case is that of a semi-positive line bundle $E$. If we let $0 \leqslant$ $\lambda_{1}(x) \leqslant \ldots \leqslant \lambda_{n}(x)$ be the eigenvalues of $\mathrm{i} \Theta(E)_{x}$ with respect to $\omega_{x}$ for all $x \in X$, formula VI-8.3 implies

$$
\begin{align*}
\left\langle A_{q} u, u\right\rangle & \geqslant\left(\lambda_{1}+\cdots+\lambda_{q}\right)|u|^{2} \\
\int_{X}\left\langle A_{q}^{-1} g, g\right\rangle d V & \leqslant \int_{X} \frac{1}{\lambda_{1}+\cdots+\lambda_{q}}|g|^{2} d V \tag{6.4}
\end{align*}
$$

A typical situation where these estimates can be applied is the case when $E$ is the trivial line bundle $X \times \mathbb{C}$ with metric given by a weight $e^{-\varphi}$. One can assume for example that $\varphi$ is plurisubharmonic and that $\mathrm{i} d^{\prime} d^{\prime \prime} \varphi$ has at least $n-q+1$ positive eigenvalues at every point, i.e. $\lambda_{q}>0$ on $X$. This situation leads to very important $L^{2}$ estimates, which are precisely those given by [Hörmander 1965, 1966]. We state here a slightly more general result.
(6.5) Theorem. Let $(X, \omega)$ be a weakly pseudoconvex Kähler manifold, $E$ a hermitian line bundle on $X, \varphi \in \mathscr{C}^{\infty}(X, \mathbb{R})$ a weight function such that the eigenvalues $\lambda_{1} \leqslant \ldots \leqslant$ $\lambda_{n}$ of $\mathrm{i} \Theta(E)+\mathrm{i} d^{\prime} d^{\prime \prime} \varphi$ are $\geqslant 0$. Then for every form $g$ of type $(n, q), q \geqslant 1$, with $L_{\mathrm{loc}}^{2}$ (resp. $\mathscr{C}^{\infty}$ ) coefficients such that $D^{\prime \prime} g=0$ and

$$
\int_{X} \frac{1}{\lambda_{1}+\cdots+\lambda_{q}}|g|^{2} e^{-\varphi} d V<+\infty
$$

we can find a $L_{\text {loc }}^{2}\left(\right.$ resp. $\left.\mathscr{C}^{\infty}\right)$ form $f$ of type $(n, q-1)$ such that $D^{\prime \prime} f=g$ and

$$
\int_{X}|f|^{2} e^{-\varphi} d V \leqslant \int_{X} \frac{1}{\lambda_{1}+\cdots+\lambda_{q}}|g|^{2} e^{-\varphi} d V
$$

Proof. Apply the general estimates to the bundle $E_{\varphi}$ deduced from $E$ by multiplication of the metric by $e^{-\varphi}$; we have $\mathrm{i} \Theta\left(E_{\varphi}\right)=\mathrm{i} \Theta(E)+\mathrm{i} d^{\prime} d^{\prime \prime} \varphi$. It is not necessary here to assume in addition that $g \in L_{n, q}^{2}\left(X, E_{\varphi}\right)$. In fact, $g$ is in $L_{\text {loc }}^{2}$ and we can exhaust X by the relatively compact weakly pseudoconvex domains

$$
X_{c}=\{x \in X ; \psi(x)<c\}
$$

where $\psi \in \mathscr{C}^{\infty}(X, \mathbb{R})$ is a plurisubharmonic exhaustion function (note that $-\log (c-\psi)$ is also such a function on $X_{c}$ ). We get therefore solutions $f_{c}$ on $X_{c}$ with uniform $L^{2}$ bounds; any weak limit $f$ gives the desired solution.

If estimates for $(p, q)$-forms instead of $(n, q)$-forms are needed, one can invoke the isomorphism $\Lambda^{p} T_{X}^{\star} \simeq \Lambda^{n-p} T_{X} \otimes \Lambda^{n} T_{X}^{\star}$ (obtained through contraction of $n$-forms by ( $n-p$ )-vectors) to get

$$
\Lambda^{p, q} T_{X}^{\star} \otimes E \simeq \Lambda^{n, q} T_{X}^{\star} \otimes F, \quad F=E \otimes \Lambda^{n-p} T_{X} .
$$

Let us look more carefully to the case $p=0$. The $(1,1)$-curvature form of $\Lambda^{n} T_{X}$ with respect to a hermitian metric $\omega$ on $T_{X}$ is called the Ricci curvature of $\omega$. We denote:
(6.6) Definition. $\operatorname{Ricci}(\omega)=\mathrm{i} \Theta\left(\Lambda^{n} T_{X}\right)=\mathrm{i} \operatorname{Tr} \Theta\left(T_{X}\right)$.

For any local coordinate system $\left(z_{1}, \ldots, z_{n}\right)$, the holomorphic $n$-form $d z_{1} \wedge \ldots \wedge d z_{n}$ is a section of $\Lambda^{n} T_{X}^{\star}$, hence Formula V-13.3 implies

$$
\begin{equation*}
\operatorname{Ricci}(\omega)=\mathrm{i} d^{\prime} d^{\prime \prime} \log \left|d z_{1} \wedge \ldots \wedge d z_{n}\right|_{\omega}^{2}=-\mathrm{i} d^{\prime} d^{\prime \prime} \log \operatorname{det}\left(\omega_{j k}\right) . \tag{6.7}
\end{equation*}
$$

The estimates of Th. 6.5 can therefore be applied to any $(0, q)$-form $g$, but $\lambda_{1} \leqslant \ldots \leqslant \lambda_{n}$ must be replaced by the eigenvalues of the ( 1,1 )-form

$$
\begin{equation*}
\mathrm{i} \Theta(E)+\operatorname{Ricci}(\omega)+\mathrm{i} d^{\prime} d^{\prime \prime} \varphi \quad(\text { supposed } \geqslant 0) \tag{6.8}
\end{equation*}
$$

We consider now domains $\Omega \subset \mathbb{C}^{n}$ equipped with the euclidean metric of $\mathbb{C}^{n}$, and the trivial bundle $E=\Omega \times \mathbb{C}$. The following result is especially convenient because it requires only weak plurisubharmonicity and avoids to compute the curvature eigenvalues.
(6.9) Theorem. Let $\Omega \subset \mathbb{C}^{n}$ be a weakly pseudoconvex open subset and $\varphi$ an upper semi-continuous plurisubharmonic function on $\Omega$. For every $\varepsilon \in] 0,1]$ and every $g \in$ $L_{p, q}^{2}(\Omega, \operatorname{loc})$ such that $d^{\prime \prime} g=0$ and

$$
\int_{\Omega}\left(1+|z|^{2}\right)|g|^{2} e^{-\varphi} d V<+\infty
$$

we can find a $L_{\text {loc }}^{2}$ form $f$ of type $(p, q-1)$ such that $d^{\prime \prime} f=g$ and

$$
\int_{\Omega}\left(1+|z|^{2}\right)^{-\varepsilon}|f|^{2} e^{-\varphi} d V \leqslant \frac{4}{q \varepsilon^{2}} \int_{\Omega}\left(1+|z|^{2}\right)|g|^{2} e^{-\varphi} d V<+\infty
$$

Moreover $f$ can be chosen smooth if $g$ and $\varphi$ are smooth.
Proof. Since $\Lambda^{p} T \Omega$ is a trivial bundle with trivial metric, the proof is immediately reduced to the case $p=0$ (or equivalently $p=n$ ). Let us first suppose that $\varphi$ is smooth. We replace $\varphi$ by $\Phi=\varphi+\tau$ where

$$
\tau(z)=\log \left(1+\left(1+|z|^{2}\right)^{\varepsilon}\right)
$$

(6.10) Lemma. The smallest eigenvalue $\lambda_{1}(z)$ of $\mathrm{i} d^{\prime} d^{\prime \prime} \tau(z)$ satisfies

$$
\lambda_{1}(z) \geqslant \frac{\varepsilon^{2}}{2\left(1+|z|^{2}\right)\left(1+\left(1+|z|^{2}\right)^{\varepsilon}\right)}
$$

In fact a brute force computation of the complex hessian $H \tau_{z}(\xi)$ and the CauchySchwarz inequality yield

$$
\begin{aligned}
& H \tau_{z}(\xi)= \\
& =\frac{\varepsilon\left(1+|z|^{2}\right)^{\varepsilon-1}|\xi|^{2}}{1+\left(1+|z|^{2}\right)^{\varepsilon}}+\frac{\varepsilon(\varepsilon-1)\left(1+|z|^{2}\right)^{\varepsilon-2}|\langle\xi, z\rangle|^{2}}{1+\left(1+|z|^{2}\right)^{\varepsilon}}-\frac{\varepsilon^{2}\left(1+|z|^{2}\right)^{2 \varepsilon-2}|\langle\xi, z\rangle|^{2}}{\left(1+\left(1+|z|^{2}\right)^{\varepsilon}\right)^{2}} \\
& \geqslant \varepsilon\left(\frac{\left(1+|z|^{2}\right)^{\varepsilon-1}}{1+\left(1+|z|^{2}\right)^{\varepsilon}}-\frac{(1-\varepsilon)\left(1+|z|^{2}\right)^{\varepsilon-2}|z|^{2}}{1+\left(1+|z|^{2}\right)^{\varepsilon}}-\frac{\varepsilon\left(1+|z|^{2}\right)^{2 \varepsilon-2}|z|^{2}}{\left(1+\left(1+|z|^{2}\right)^{\varepsilon}\right)^{2}}\right)|\xi|^{2} \\
& =\varepsilon \frac{1+\varepsilon|z|^{2}+\left(1+|z|^{2}\right)^{\varepsilon}}{\left(1+|z|^{2}\right)^{2-\varepsilon}\left(1+\left(1+|z|^{2}\right)^{\varepsilon}\right)^{2}}|\xi|^{2} \geqslant \frac{\varepsilon^{2}|\xi|^{2}}{\left(1+|z|^{2}\right)^{1-\varepsilon}\left(1+\left(1+|z|^{2}\right)^{\varepsilon}\right)^{2}} \\
& \geqslant \frac{\varepsilon^{2}}{2\left(1+|z|^{2}\right)\left(1+\left(1+|z|^{2}\right)^{\varepsilon}\right)}|\xi|^{2} .
\end{aligned}
$$

The Lemma implies $e^{-\tau} / \lambda_{1} \leqslant 2\left(1+|z|^{2}\right) / \varepsilon^{2}$, thus Cor. 6.5 provides an $f$ such that

$$
\int_{\Omega}\left(1+\left(1+|z|^{2}\right)^{\varepsilon}\right)^{-1}|f|^{2} e^{-\varphi} d V \leqslant \frac{2}{q \varepsilon^{2}} \int_{\Omega}\left(1+|z|^{2}\right)|g|^{2} e^{-\varphi} d V<+\infty
$$

and the required estimate follows. If $\varphi$ is not smooth, apply the result to a sequence of regularized weights $\rho_{\varepsilon} \star \varphi \geqslant \varphi$ on an increasing sequence of domains $\Omega_{c} \subset \subset \Omega$, and extract a weakly convergent subsequence of solutions.

## § 7. Extension of Holomorphic Functions from Subvarieties

The existence theorems for solutions of the $d^{\prime \prime}$ operator easily lead to an extension theorem for sections of a holomorphic line bundle defined in a neighborhood of an analytic subset. The following result [Demailly 1982] is an improvement and a generalization of Jennane's extension theorem [Jennane 1976].
(7.1) Theorem. Let $(X, \omega)$ be a weakly pseudoconvex Kähler manifold, $L$ a hermitian line bundle and $E$ a hermitian vector bundle over $X$. Let $Y$ be an analytic subset of $X$ such that $Y=\sigma^{-1}(0)$ for some section $\sigma$ of $E$, and $p$ the maximal codimension of the irreducible components of $Y$. Let $f$ be a holomorphic section of $K_{X} \otimes L$ defined in the open set $U \supset Y$ of points $x \in X$ such that $|\sigma(x)|<1$. If $\int_{U}|f|^{2} d V<+\infty$ and if the curvature form of $L$ satisfies

$$
\mathrm{i} \Theta(L) \geqslant\left(\frac{p}{|\sigma|^{2}}+\frac{\varepsilon}{1+|\sigma|^{2}}\right)\{\mathrm{i} \Theta(E) \sigma, \sigma\}
$$

for some $\varepsilon>0$, there is a section $F \in H^{0}\left(X, K_{X} \otimes L\right)$ such that $F_{\uparrow Y}=f_{\uparrow Y}$ and

$$
\int_{X} \frac{|F|^{2}}{\left(1+|\sigma|^{2}\right)^{p+\varepsilon}} d V \leqslant\left(1+\frac{(p+1)^{2}}{\varepsilon}\right) \int_{U}|f|^{2} d V
$$

The proof will involve a weight with logarithmic singularities along $Y$. We must therefore apply the existence theorem over $X \backslash Y$. This requires to know whether $X \backslash Y$ has a complete Kähler metric.
(7.2) Lemma. Let $(X, \omega)$ be a Kähler manifold, and $Y=\sigma^{-1}(0)$ an analytic subset defined by a section of a hermitian vector bundle $E \rightarrow X$. If $X$ is weakly pseudoconvex and exhausted by $X_{c}=\{x \in X ; \psi(x)<c\}$, then $X_{c} \backslash Y$ has a complete Kähler metric for all $c \in \mathbb{R}$. The same conclusion holds for $X \backslash Y$ if $(X, \omega)$ is complete and if for some constant $C \geqslant 0$ we have $\Theta_{E} \leqslant_{\text {Grif }} C \omega \otimes\langle,\rangle_{E}$ on $X$.

Proof. Set $\tau=\log |\sigma|^{2}$. Then $d^{\prime} \tau=\left\{D^{\prime} \sigma, \sigma\right\} /|\sigma|^{2}$ and $D^{\prime \prime} D^{\prime} \sigma=D^{2} \sigma=\Theta(E) \sigma$, thus

$$
\mathrm{i} d^{\prime} d^{\prime \prime} \tau=\mathrm{i} \frac{\left\{D^{\prime} \sigma, D^{\prime} \sigma\right\}}{|\sigma|^{2}}-\mathrm{i} \frac{\left\{D^{\prime} \sigma, \sigma\right\} \wedge\left\{\sigma, D^{\prime} \sigma\right\}}{|\sigma|^{4}}-\frac{\{\mathrm{i} \Theta(E) \sigma, \sigma\}}{|\sigma|^{2}} .
$$

For every $\xi \in T_{X}$, we find therefore

$$
\begin{aligned}
H \tau(\xi) & =\frac{|\sigma|^{2}\left|D^{\prime} \sigma \cdot \xi\right|^{2}-\left|\left\langle D^{\prime} \sigma \cdot \xi, \sigma\right\rangle\right|^{2}}{|\sigma|^{4}}-\frac{\Theta_{E}(\xi \otimes \sigma, \xi \otimes \sigma)}{|\sigma|^{2}} \\
& \geqslant-\frac{\Theta_{E}(\xi \otimes \sigma, \xi \otimes \sigma)}{|\sigma|^{2}}
\end{aligned}
$$

by the Cauchy-Schwarz inequality. If $C$ is a bound for the coefficients of $\Theta_{E}$ on the compact subset $\bar{X}_{c}$, we get $\mathrm{i} d^{\prime} d^{\prime \prime} \tau \geqslant-C \omega$ on $X_{c}$. Let $\chi \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$ be a convex increasing function. We set

$$
\widehat{\omega}=\omega+\mathrm{id} d^{\prime} d^{\prime \prime}(\chi \circ \tau) .
$$

Formula 5.3 shows that $\widehat{\omega}$ is positive definite if $\chi^{\prime} \leqslant 1 / 2 C$ and that $\widehat{\omega}$ is complete near $Y=\tau^{-1}(-\infty)$ as soon as

$$
\int_{-\infty}^{0} \sqrt{\chi^{\prime \prime}(t)} d t=+\infty
$$

One can choose for example $\chi$ such that $\chi(t)=\frac{1}{5 C}(t-\log |t|)$ for $t \leqslant-1$. In order to obtain a complete Kähler metric on $X_{c} \backslash Y$, we need also that the metric be complete near $\partial X_{c}$. Such a metric is given by

$$
\begin{aligned}
\widetilde{\omega} & =\widehat{\omega}+\mathrm{i} d^{\prime} d^{\prime \prime} \log (c-\psi)^{-1}=\widehat{\omega}+\frac{\mathrm{i} d^{\prime} d^{\prime \prime} \psi}{c-\psi}+\frac{\mathrm{i} d^{\prime} \psi \wedge d^{\prime \prime} \psi}{(c-\psi)^{2}} \\
& \geqslant \mathrm{i} d^{\prime} \log (c-\psi)^{-1} \wedge d^{\prime \prime} \log (c-\psi)^{-1}
\end{aligned}
$$

$\widetilde{\omega}$ is complete on $X_{c} \backslash \Omega$ because $\log (c-\psi)^{-1}$ tends to $+\infty$ on $\partial X_{c}$.
Proof of Theorem 7.1.. When we replace $\sigma$ by $(1+\eta) \sigma$ for some small $\eta>0$ and let $\eta$ tend to 0 , we see that we can assume $f$ defined in a neighborhood of $\bar{U}$. Let $h$ be the continuous section of $L$ such that $h=\left(1-|\sigma|^{p+1}\right) f$ on $U=\{|\sigma|<1\}$ and $h=0$ on $X \backslash U$. We have $h_{\upharpoonright Y}=f_{\uparrow Y}$ and

$$
d^{\prime \prime} h=-\frac{p+1}{2}|\sigma|^{p-1}\left\{\sigma, D^{\prime} \sigma\right\} f \quad \text { on } \quad U, \quad d^{\prime \prime} h=0 \quad \text { on } \quad X \backslash U .
$$

We consider $g=d^{\prime \prime} h$ as a $(n, 1)$-form with values in the hermitian line bundle $L_{\varphi}=L$, endowed with the weight $e^{-\varphi}$ given by

$$
\varphi=p \log |\sigma|^{2}+\varepsilon \log \left(1+|\sigma|^{2}\right)
$$

Notice that $\varphi$ is singular along $Y$. The Cauchy-Schwarz inequality implies i $\left\{D^{\prime} \sigma, \sigma\right\} \wedge$ $\left\{\sigma, D^{\prime} \sigma\right\} \leqslant \mathrm{i}\left\{D^{\prime} \sigma, D^{\prime} \sigma\right\}$ as in Lemma 7.2, and we find

$$
\begin{aligned}
\mathrm{i} d^{\prime} d^{\prime \prime} \log \left(1+|\sigma|^{2}\right)= & \frac{\left(1+|\sigma|^{2}\right) \mathrm{i}\left\{D^{\prime} \sigma, D^{\prime} \sigma\right\}-\mathrm{i}\left\{D^{\prime} \sigma, \sigma\right\} \wedge\left\{\sigma, D^{\prime} \sigma\right\}}{\left(1+|\sigma|^{2}\right)^{2}} \\
& -\frac{\{\mathrm{i} \Theta(E) \sigma, \sigma\}}{1+|\sigma|^{2}} \geqslant \frac{\mathrm{i}\left\{D^{\prime} \sigma, D^{\prime} \sigma\right\}}{\left(1+|\sigma|^{2}\right)^{2}}-\frac{\{\mathrm{i} \Theta(E) \sigma, \sigma\}}{1+|\sigma|^{2}}
\end{aligned}
$$

The inequality i $d^{\prime} d^{\prime \prime} \log |\sigma|^{2} \geqslant-\{\mathrm{i} \Theta(E) \sigma, \sigma\} /|\sigma|^{2}$ obtained in Lemma 7.2 and the above one imply

$$
\begin{aligned}
\mathrm{i} \Theta\left(L_{\varphi}\right) & =\mathrm{i} \Theta(L)+p \mathrm{i} d^{\prime} d^{\prime \prime} \log |\sigma|^{2}+\varepsilon \mathrm{i} d^{\prime} d^{\prime \prime} \log \left(1+|\sigma|^{2}\right) \\
& \geqslant \mathrm{i} \Theta(L)-\left(\frac{p}{|\sigma|^{2}}+\frac{\varepsilon}{1+|\sigma|^{2}}\right)\{\mathrm{i} \Theta(E) \sigma, \sigma\}+\varepsilon \frac{\mathrm{i}\left\{D^{\prime} \sigma, D^{\prime} \sigma\right\}}{\left(1+|\sigma|^{2}\right)^{2}} \\
& \geqslant \varepsilon \frac{\mathrm{i}\left\{D^{\prime} \sigma, \sigma\right\} \wedge\left\{\sigma, D^{\prime} \sigma\right\}}{|\sigma|^{2}\left(1+|\sigma|^{2}\right)^{2}},
\end{aligned}
$$

thanks to the hypothesis on the curvature of $L$ and the Cauchy-Schwarz inequality. Set $\xi=(p+1) / 2|\sigma|^{p-1}\left\{D^{\prime} \sigma, \sigma\right\}=\sum \xi_{j} d z_{j}$ in an $\omega$-orthonormal basis $\partial / \partial z_{j}$, and let
$\widehat{\xi}=\sum \xi_{j} \partial / \partial \bar{z}_{j}$ be the dual $(0,1)$-vector field. For every $(n, 1)$-form $v$ with values in $L_{\varphi}$, we find

$$
\begin{aligned}
\left|\left\langle d^{\prime \prime} h, v\right\rangle\right| & =|\langle\bar{\xi} \wedge f, v\rangle|=|\langle f, \widehat{\xi}\lrcorner v\rangle|\leqslant|f|| \widehat{\xi}\lrcorner v \mid, \\
\widehat{\xi}\lrcorner v & =\sum-\mathrm{i} \xi_{j} d z_{j} \wedge \Lambda v=-\mathrm{i} \xi \wedge \Lambda v, \\
\left|\left\langle d^{\prime \prime} h, v\right\rangle\right|^{2} & \left.\left.\leqslant|f|^{2} \mid \widehat{\xi}\right\lrcorner\left.v\right|^{2}=|f|^{2}\langle-\mathrm{i} \xi \wedge \Lambda v, \widehat{\xi}\lrcorner v\right\rangle \\
& =|f|^{2}\langle-\mathrm{i} \bar{\xi} \wedge \xi \wedge \Lambda v, v\rangle=|f|^{2}\langle[\mathrm{i} \xi \wedge \bar{\xi}, \Lambda] v, v\rangle \\
& \leqslant \frac{(p+1)^{2}}{4 \varepsilon}|\sigma|^{2 p}\left(1+|\sigma|^{2}\right)^{2}|f|^{2}\left\langle\left[\mathrm{i} \Theta\left(L_{\varphi}\right), \Lambda\right] v, v\right\rangle .
\end{aligned}
$$

Thus, in the notations of Th. 6.1, the form $g=d^{\prime \prime} h$ satisfies

$$
\left\langle A_{1}^{-1} g, g\right\rangle \leqslant \frac{(p+1)^{2}}{4 \varepsilon}|\sigma|^{2 p}\left(1+|\sigma|^{2}\right)^{2}|f|^{2} \leqslant \frac{(p+1)^{2}}{\varepsilon}|f|^{2} e^{\varphi},
$$

where the last equality results from the fact that $\left(1+|\sigma|^{2}\right)^{2} \leqslant 4$ on the support of $g$. Lemma 7.2 shows that the existence theorem 6.1 can be applied on each set $X_{c} \backslash Y$. Letting $c$ tend to infinity, we infer the existence of a $(n, 0)$-form $u$ with values in $L$ such that $d^{\prime \prime} u=g$ on $X \backslash Y$ and

$$
\begin{aligned}
\int_{X \backslash Y}|u|^{2} e^{-\varphi} d V & \leqslant \int_{X \backslash Y}\left\langle A_{1}^{-1} g, g\right\rangle e^{-\varphi}, \quad \text { thus } \\
\int_{X \backslash Y} \frac{|u|^{2}}{|\sigma|^{2 p}\left(1+|\sigma|^{2}\right)^{\varepsilon}} d V & \leqslant \frac{(p+1)^{2}}{\varepsilon} \int_{U}|f|^{2} d V .
\end{aligned}
$$

This estimate implies in particular that $u$ is locally $L^{2}$ near $Y$. As $g$ is continuous over $X$, Lemma 7.3 below shows that the equality $d^{\prime \prime} u=g=d^{\prime \prime} h$ extends to $X$, thus $F=h-u$ is holomorphic everywhere. Hence $u=h-F$ is continuous on $X$. As $|\sigma(x)| \leqslant C d(x, Y)$ in a neighborhood of every point of $Y$, we see that $|\sigma|^{-2 p}$ is non integrable at every point $x_{0} \in Y_{\text {reg }}$, because codim $Y \leqslant p$. It follows that $u=0$ on $Y$, so

$$
F_{\upharpoonright Y}=h_{\upharpoonright Y}=f_{\upharpoonright Y} .
$$

The final $L^{2}$-estimate of Th. 7.1 follows from the inequality

$$
|F|^{2}=|h-u|^{2} \leqslant\left(1+|\sigma|^{-2 p}\right)|u|^{2}+\left(1+|\sigma|^{2 p}\right)|f|^{2}
$$

which implies

$$
\frac{|F|^{2}}{\left(1+|\sigma|^{2}\right)^{p}} \leqslant \frac{|u|^{2}}{|\sigma|^{2 p}}+|f|^{2} .
$$

(7.3) Lemma. Let $\Omega$ be an open subset of $\mathbb{C}^{n}$ and $Y$ an analytic subset of $\Omega$. Assume that $v$ is a $(p, q-1)$-form with $L_{\mathrm{loc}}^{2}$ coefficients and $w a(p, q)$-form with $L_{\mathrm{loc}}^{1}$ coefficients such that $d^{\prime \prime} v=w$ on $\Omega \backslash Y$ (in the sense of distribution theory). Then $d^{\prime \prime} v=w$ on $\Omega$.

Proof. An induction on the dimension of $Y$ shows that it is sufficient to prove the result in a neighborhood of a regular point $a \in Y$. By using a local analytic isomorphism, the proof is reduced to the case where $Y$ is contained in the hyperplane $z_{1}=0$, with $a=0$.

Let $\lambda \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$ be a function such that $\lambda(t)=0$ for $t \leqslant \frac{1}{2}$ and $\lambda(t)=1$ for $t \geqslant 1$. We must show that

$$
\begin{equation*}
\int_{\Omega} w \wedge \alpha=(-1)^{p+q} \int_{\Omega} v \wedge d^{\prime \prime} \alpha \tag{7.4}
\end{equation*}
$$

for all $\alpha \in \mathscr{D}_{n-p, n-q}(\Omega)$. Set $\lambda_{\varepsilon}(z)=\lambda\left(\left|z_{1}\right| / \varepsilon\right)$ and replace $\alpha$ in the integral by $\lambda_{\varepsilon} \alpha$. Then $\lambda_{\varepsilon} \alpha \in \mathscr{D}_{n-p, n-q}(\Omega \backslash Y)$ and the hypotheses imply

$$
\int_{\Omega} w \wedge \lambda_{\varepsilon} \alpha=(-1)^{p+q} \int_{\Omega} v \wedge d^{\prime \prime}\left(\lambda_{\varepsilon} \alpha\right)=(-1)^{p+q} \int_{\Omega} v \wedge\left(d^{\prime \prime} \lambda_{\varepsilon} \wedge \alpha+\lambda_{\varepsilon} d^{\prime \prime} \alpha\right) .
$$

As $w$ and $v$ have $L_{\text {loc }}^{1}$ coefficients on $\Omega$, the integrals of $w \wedge \lambda_{\varepsilon} \alpha$ and $v \wedge \lambda_{\varepsilon} d^{\prime \prime} \alpha$ converge respectively to the integrals of $w \wedge \alpha$ and $v \wedge d^{\prime \prime} \alpha$ as $\varepsilon$ tends to 0 . The remaining term can be estimated by means of the Cauchy-Schwarz inequality:

$$
\left|\int_{\Omega} v \wedge d^{\prime \prime} \lambda_{\varepsilon} \wedge \alpha\right|^{2} \leqslant \int_{\left|z_{1}\right| \leqslant \varepsilon}|v \wedge \alpha|^{2} d V \cdot \int_{\operatorname{Supp} \alpha}\left|d^{\prime \prime} \lambda_{\varepsilon}\right|^{2} d V
$$

as $v \in L_{\text {loc }}^{2}(\Omega)$, the integral $\int_{\left|z_{1}\right| \leqslant \varepsilon}|v \wedge \alpha|^{2} d V$ converges to 0 with $\varepsilon$, whereas

$$
\int_{\operatorname{Supp} \alpha}\left|d^{\prime \prime} \lambda_{\varepsilon}\right|^{2} d V \leqslant \frac{C}{\varepsilon^{2}} \operatorname{Vol}\left(\operatorname{Supp} \alpha \cap\left\{\left|z_{1}\right| \leqslant \varepsilon\right\}\right) \leqslant C^{\prime}
$$

Equality (7.4) follows when $\varepsilon$ tends to 0 .
(7.5) Corollary. Let $\Omega \subset \mathbb{C}^{n}$ be a weakly pseudoconvex domain and let $\varphi, \psi$ be plurisubharmonic functions on $\Omega$, where $\psi$ is supposed to be finite and continuous. Let $\sigma=\left(\sigma_{1}, \ldots, \sigma_{r}\right)$ be a family of holomorphic functions on $\Omega$, let $Y=\sigma^{-1}(0), p=$ maximal codimension of $Y$ and set
a) $U=\left\{z \in \Omega ;|\sigma(z)|^{2}<e^{-\psi(z)}\right\}$, resp.
b) $U^{\prime}=\left\{z \in \Omega ;|\sigma(z)|^{2}<e^{\psi(z)}\right\}$.

For every $\varepsilon>0$ and every holomorphic function $f$ on $U$, there exists a holomorphic function $F$ on $\Omega$ such that $F_{\upharpoonright Y}=f_{\lceil Y}$ and
a) $\quad \int_{\Omega} \frac{|F|^{2} e^{-\varphi+p \psi}}{\left(1+|\sigma|^{2} e^{\psi}\right)^{p+\varepsilon}} d V \leqslant\left(1+\frac{(p+1)^{2}}{\varepsilon}\right) \int_{U}|f|^{2} e^{-\varphi+p \psi} d V, \quad$ resp.
b) $\quad \int_{\Omega} \frac{|F|^{2} e^{-\varphi}}{\left(e^{\psi}+|\sigma|^{2}\right)^{p+\varepsilon}} d V \leqslant\left(1+\frac{(p+1)^{2}}{\varepsilon}\right) \int_{U}|f|^{2} e^{-\varphi-(p+\varepsilon) \psi} d V$.

Proof. After taking convolutions with smooth kernels on pseudoconvex subdomains $\Omega_{c} \subset$ $\subset \Omega$, we may assume $\varphi, \psi$ smooth. In either case a) or b), apply Th. 7.1 to
a) $E=\Omega \times \mathbb{C}^{r}$ with the weight $e^{\psi}, L=\Omega \times \mathbb{C}$ with the weight $e^{-\varphi+p \psi}$, and $U=$ $\left\{|\sigma|^{2} e^{\psi}<1\right\}$. Then

$$
\mathrm{i} \Theta(E)=-\mathrm{i} d^{\prime} d^{\prime \prime} \psi \otimes \operatorname{Id}_{E} \leqslant 0, \quad \mathrm{i} \Theta(L)=\mathrm{i} d^{\prime} d^{\prime \prime} \varphi-p \mathrm{i} d^{\prime} d^{\prime \prime} \psi \geqslant p \mathrm{i} \Theta(E)
$$

b) $E=\Omega \times \mathbb{C}^{r}$ with the weight $e^{-\psi}, L=\Omega \times \mathbb{C}$ with the weight $e^{-\varphi-(p+\varepsilon) \psi}$, and $U=\left\{|\sigma|^{2} e^{-\psi}<1\right\}$. Then

$$
\mathrm{i} \Theta(E)=\mathrm{i} d^{\prime} d^{\prime \prime} \psi \otimes \operatorname{Id}_{E} \geqslant 0, \quad \mathrm{i} \Theta(L)=\mathrm{i} d^{\prime} d^{\prime \prime} \varphi+(p+\varepsilon) \mathrm{id}^{\prime} d^{\prime \prime} \psi \geqslant(p+\varepsilon) \mathrm{i} \Theta(E)
$$

The condition on $\Theta(L)$ is satisfied in both cases and $K_{\Omega}$ is trivial.
(7.6) Hörmander-Bombieri-Skoda theorem. Let $\Omega \subset \mathbb{C}^{n}$ be a weakly pseudoconvex domain and $\varphi$ a plurisubharmonic function on $\Omega$. For every $\varepsilon>0$ and every point $z_{0} \in \Omega$ such that $e^{-\varphi}$ is integrable in a neighborhood of $z_{0}$, there exists a holomorphic function $F$ on $\Omega$ such that $F\left(z_{0}\right)=1$ and

$$
\int_{\Omega} \frac{|F(z)|^{2} e^{-\varphi(z)}}{\left(1+|z|^{2}\right)^{n+\varepsilon}} d V<+\infty
$$

[Bombieri 1970] originally stated the theorem with the exponent $3 n$ instead of $n+\varepsilon$; the improved exponent $n+\varepsilon$ is due to [Skoda 1975]. The example $\Omega=\mathbb{C}^{n}, \varphi(z)=0$ shows that one cannot replace $\varepsilon$ by 0 .

Proof. Apply Cor. 7.5 b) to $f \equiv 1, \sigma(z)=z-z_{0}, p=n$ and $\psi \equiv \log r^{2}$ where $U=B\left(z_{0}, r\right)$ is a ball such that $\int_{U} e^{-\varphi} d V<+\infty$.
(7.7) Corollary. Let $\varphi$ be a plurisubharmonic function on a complex manifold $X$. Let $A$ be the set of points $z \in X$ such that $e^{-\varphi}$ is not locally integrable in a neighborhood of $z$. Then $A$ is an analytic subset of $X$.

Proof. Let $\Omega \subset X$ be an open coordinate patch isomorphic to a ball of $\mathbb{C}^{n}$, with coordinates $\left(z_{1}, \ldots, z_{n}\right)$. Define $E \subset H^{0}(\Omega, \mathscr{O})$ to be the Hilbert space of holomorphic functions $f$ on $\Omega$ such that

$$
\int_{\Omega}|f(z)|^{2} e^{-\varphi(z)} d V(z)<+\infty
$$

Then $A \cap \Omega=\bigcap_{f \in E} f^{-1}(0)$. In fact, every $f$ in $E$ must obviously vanish on $A$; conversely, if $z_{0} \notin A$, Th. 7.6 shows that there exists $f \in E$ such that $f\left(z_{0}\right) \neq 0$. By Th. II-5.5, we conclude that $A$ is analytic.

## § 8. Applications to Hypersurface Singularities

We first give some basic definitions and results concerning multiplicities of divisors on a complex manifold.
(8.1) Proposition. Let $X$ be a complex manifold and $\Delta=\sum \lambda_{j}\left[Z_{j}\right]$ a divisor on $X$ with real coefficients $\lambda_{j} \geqslant 0$. Let $x \in X$ and $f_{j}=0,1 \leqslant j \leqslant N$, irreducible equations of $Z_{j}$ on a neighborhood $U$ of $x$.
a) The multiplicity of $\Delta$ at $x$ is defined by

$$
\mu(\Delta, x)=\sum \lambda_{j} \operatorname{ord}_{x} f_{j}
$$

b) $\Delta$ is said to have normal crossings at a point $x \in \operatorname{Supp} \Delta$ if all hypersurfaces $Z_{j}$ containing $x$ are smooth at $x$ and intersect transversally, i.e. if the linear forms $d f_{j}$ defining the corresponding tangent spaces $T_{x} Z_{j}$ are linearly independent at $x$. The set $\operatorname{nnc}(\Delta)$ of non normal crossing points is an analytic subset of $X$.
c) The non-integrability locus nil( $\Delta$ ) is defined as the set of points $x \in X$ such that $\Pi\left|f_{j}\right|^{-2 \lambda_{j}}$ is non integrable near $x$. Then $\operatorname{nil}(\Delta)$ is an analytic subset of $X$ and there are inclusions

$$
\{x \in X ; \mu(\Delta, x) \geqslant n\} \subset \operatorname{nil}(\Delta) \subset\{x \in X ; \mu(\Delta, x) \geqslant 1\}
$$

Moreover $\operatorname{nil}(\Delta) \subset \operatorname{nnc}(\Delta)$ if all coefficients of $\Delta$ satisfy $\lambda_{j}<1$.
Proof. b) The set $\operatorname{nnc}(\Delta) \cap U$ is the union of the analytic sets

$$
f_{j_{1}}=\ldots=f_{j_{p}}=0, \quad d f_{j_{1}} \wedge \ldots \wedge d f_{j_{p}}=0
$$

for each subset $\left\{j_{1}, \ldots, j_{p}\right\}$ of the index set $\{1, \ldots, N\}$. Thus nnc $(\Delta)$ is analytic.
c) The analyticity of nil $(\Delta)$ follows from Cor. 7.7 applied to the plurisubharmonic function $\varphi=\sum 2 \lambda_{j} \log \left|f_{j}\right|$. Assume first that $\lambda_{j}<1$ and that $\Delta$ has normal crossings at $x$. Let $f_{j_{1}}(x)=\ldots=f_{j_{s}}(x)=0$ and $f_{j}(x) \neq 0$ for $j \neq j_{l}$. Then, we can choose local coordinates $\left(w_{1}, \ldots, w_{n}\right)$ on $U$ such that $w_{1}=f_{j_{1}}(z), \ldots, w_{s}=f_{j_{s}}(z)$, and we have

$$
\int_{U} \frac{d \lambda(z)}{\prod\left|f_{j}(z)\right|^{2 \lambda_{j}}} \leqslant \int_{U} \frac{C d \lambda(w)}{\left|w_{1}\right|^{2 \lambda_{1}} \ldots\left|w_{s}\right|^{2 \lambda_{s}}}<+\infty
$$

It follows that $\operatorname{nil}(\Delta) \subset \operatorname{nnc}(\Delta)$. Let us prove now the statement relating nil( $\Delta$ ) with multiplicity sets. Near any point $x$, we have $\left|f_{j}(z)\right| \leqslant C_{j}|z-x|^{m_{j}}$ with $m_{j}=\operatorname{ord}_{x} f_{j}$, thus

$$
\prod\left|f_{j}\right|^{-2 \lambda_{j}} \geqslant C|z-x|^{-2 \mu(\Delta, x)}
$$

It follows that $x \in \operatorname{nil}(\Delta)$ as soon as $\mu(\Delta, x) \geqslant n$. On the other hand, we are going to prove that $\mu(\Delta, x)<1$ implies $x \notin \operatorname{nil}(\Delta)$, i.e. $\prod\left|f_{j}\right|^{-2 \lambda_{j}}$ integrable near $x$. We may assume $\lambda_{j}$ rational; otherwise replace each $\lambda_{j}$ by a slightly larger rational number in such a way that $\mu(\Delta, x)<1$ is still true. Set $f=\prod f_{j}^{k \lambda_{j}}$ where $k$ is a common denominator. The result is then a consequence of the following lemma.
(8.2) Lemma. If $f \in \mathscr{O}_{X, x}$ is not identically 0 , there exists a neighborhood $U$ of $x$ such that $\int_{U}|f|^{-2 \lambda} d V$ converges for all $\lambda<1 / m, m=\operatorname{ord}_{x} f$.

Proof. One can assume that $f$ is a Weierstrass polynomial

$$
f(z)=z_{n}^{m}+a_{1}\left(z^{\prime}\right) z_{n}^{m-1}+\cdots+a_{m}\left(z^{\prime}\right), \quad a_{j}\left(z^{\prime}\right) \in \mathbb{O}_{n-1}, \quad a_{j}(0)=0,
$$

with respect to some coordinates $\left(z_{1}, \ldots, z_{n}\right)$ centered at $x$. Let $v_{j}\left(z^{\prime}\right), 1 \leqslant j \leqslant m$, denote the roots $z_{n}$ of $f(z)=0$. On a small neighborhood $U$ of $x$ we have $\left|v_{j}\left(z^{\prime}\right)\right| \leqslant 1$. The inequality between arithmetic and geometric mean implies

$$
\begin{aligned}
\int_{\left\{\left|z_{n}\right| \leqslant 1\right\}}|f(z)|^{-2 \lambda} d x_{n} d y_{n} & =\int_{\left\{\left|z_{n}\right| \leqslant 1\right\}} \prod_{1 \leqslant j \leqslant m}\left|z_{n}-v_{j}\left(z^{\prime}\right)\right|^{-2 \lambda} d x_{n} d y_{n} \\
& \leqslant \frac{1}{m} \int_{\left\{\left|z_{n}\right| \leqslant 1\right\}} \sum_{1 \leqslant j \leqslant m}\left|z_{n}-v_{j}\left(z^{\prime}\right)\right|^{-2 m \lambda} d x_{n} d y_{n} \\
& \leqslant \int_{\left\{\left|z_{n}\right| \leqslant 2\right\}} \frac{d x_{n} d y_{n}}{\left|z_{n}\right|^{2 m \lambda}}
\end{aligned}
$$

so the Lemma follows from the Fubini theorem.
Another interesting application concerns the study of multiplicities of singular points for algebraic hypersurfaces in $\mathbb{P}^{n}$. Following [Waldschmidt 1975], we introduce the following definition.
(8.3) Definition. Let $S$ be a finite subset of $\mathbb{P}^{n}$. For any integer $t \geqslant 1$, we define $\omega_{t}(S)$ as the minimum of the degrees of non zero homogeneous polynomials $P \in \mathbb{C}\left[z_{0}, \ldots, z_{n}\right]$ which vanish at order $t$ at every point of S, i.e. $D^{\alpha} P(w)=0$ for every $w \in S$ and every multi-index $\alpha=\left(\alpha_{0}, \ldots, \alpha_{n}\right)$ of length $|\alpha|<t$.

It is clear that $t \longmapsto \omega_{t}(S)$ is a non-decreasing and subadditive function, i.e. for all integers $t_{1}, t_{2} \geqslant 1$ we have $\omega_{t_{1}+t_{2}}(S) \leqslant \omega_{t_{1}}(S)+\omega_{t_{2}}(S)$. One defines

$$
\begin{equation*}
\Omega(S)=\inf _{t \geqslant 1} \frac{\omega_{t}(S)}{t} \tag{8.4}
\end{equation*}
$$

For all integers $t, t^{\prime} \geqslant 1$, the monotonicity and subadditivity of $\omega_{t}(S)$ show that

$$
\omega_{t}(S) \leqslant\left(\left[t / t^{\prime}\right]+1\right) \omega_{t^{\prime}}(S), \quad \text { hence } \Omega(S) \leqslant \frac{\omega_{t}(S)}{t} \leqslant\left(\frac{1}{t^{\prime}}+\frac{1}{t}\right) \omega_{t^{\prime}}(S)
$$

We find therefore

$$
\begin{equation*}
\Omega(S)=\lim _{t \rightarrow+\infty} \frac{\omega_{t}(S)}{t} \tag{8.5}
\end{equation*}
$$

Our goal is to find a lower bound of $\Omega(S)$ in terms of $\omega_{t}(S)$. For $n=1$, it is obvious that $\Omega(S)=\omega_{t}(S) / t=\operatorname{card} S$ for all $t$. From now on, we assume that $n \geqslant 2$.
(8.6) Theorem. Let $t_{1}, t_{2} \geqslant 1$ be integers, let $P$ be a homogeneous polynomial of degree $\omega_{t_{2}}(S)$ vanishing at order $\geqslant t_{2}$ at every point of $S$. If $P=P_{1}^{k_{1}} \ldots P_{N}^{k_{N}}$ is the decomposition of $P$ in irreducible factors and $Z_{j}=P_{j}^{-1}(0)$, we set

$$
\alpha=\frac{t_{1}+n-1}{t_{2}}, \quad \Delta=\sum\left(k_{j} \alpha-\left[k_{j} \alpha\right]\right)\left[Z_{j}\right], \quad a=\operatorname{dim}(\operatorname{nil}(\Delta)) .
$$

Then we have the inequality

$$
\frac{\omega_{t_{1}}(S)+n-a-1}{t_{1}+n-1} \leqslant \frac{\omega_{t_{2}}(S)}{t_{2}} .
$$

Let us first make a few comments before giving the proof. If we let $t_{2}$ tend to infinity and observe that $\operatorname{nil}(\Delta) \subset \operatorname{nnc}(\Delta)$ by Prop. 8.1 c ), we get $a \leqslant 2$ and

$$
\begin{equation*}
\frac{\omega_{t_{1}}(S)+1}{t_{1}+n-1} \leqslant \Omega(S) \leqslant \frac{\omega_{t_{2}}(S)}{t_{2}} \tag{8.7}
\end{equation*}
$$

Such a result was first obtained by [Waldschmidt 1975, 1979] with the lower bound $\omega_{t_{1}}(S) /\left(t_{1}+n-1\right)$, as a consequence of the Hörmander-Bombieri-Skoda theorem. The
above improved inequalities were then found by (Esnault-Viehweg 1983), who used rather deep tools of algebraic geometry. Our proof will consist in a refinement of the BombieriWaldschmidt method due to [Azhari 1990]. It has been conjectured by [Chudnovsky 1979] that $\Omega(S) \geqslant\left(\omega_{1}(S)+n-1\right) / n$. Chudnovsky's conjecture is true for $n=2$ (as shown by (8.7)); this case was first verified independently by [Chudnovsky 1979] and [Demailly 1982]. The conjecture can also be verified in case $S$ is a complete polytope, and the lower bound of the conjecture is then optimal (see [Demailly 1982a] and ??.?.?). More generally, it is natural to ask whether the inequality

$$
\begin{equation*}
\frac{\omega_{t_{1}}(S)+n-1}{t_{1}+n-1} \leqslant \Omega(S) \leqslant \frac{\omega_{t_{2}}(S)}{t_{2}} \tag{8.8}
\end{equation*}
$$

always holds; this is the case if there are infinitely many $t_{2}$ for which $P$ can be chosen in such a way that $\operatorname{nil}(\Delta)$ has dimension $a=0$.
(8.9) Bertini's lemma. If $E \subset \mathbb{P}^{n}$ is an analytic subset of dimension a, there exists a dense subset in the grassmannian of $k$-codimensional linear subspaces $Y$ of $\mathbb{P}^{n}$ such that $\operatorname{dim}(E \cap Y) \leqslant a-k$ (when $k>a$ this means that $E \cap Y=\emptyset)$.

Proof. By induction on $n$, it suffices to show that $\operatorname{dim}(E \cap H) \leqslant a-1$ for a generic hyperplane $H \subset \mathbb{P}^{n}$. Let $E_{j}$ be the (finite) family of irreducible components of $E$, and $w_{j} \in E_{j}$ an arbitrary point. Then $E \cap H=\bigcup E_{j} \cap H$ and we have $\operatorname{dim} E_{j} \cap H<\operatorname{dim} E_{j} \leqslant$ $a$ as soon as $H$ avoids all points $w_{j}$.

Proof of Theorem 8.6.. By Bertini's lemma, there exists a linear subspace $Y \subset \mathbb{P}^{n}$ of codimension $a+1$ such that $\operatorname{nil}(\Delta) \cap Y=\emptyset$. We consider $P$ as a section of the line bundle $\mathscr{O}(D)$ over $\mathbb{P}^{n}$, where $D=\operatorname{deg} P$ (cf. Th. V-15.5). There are sections $\sigma_{1}, \ldots, \sigma_{a+1}$ of $\mathfrak{O}(1)$ such that $Y=\sigma^{-1}(0)$. We shall apply Th. 7.1 to $E=\mathbb{O}(1)$ with its standard hermitian metric, and to $L=\mathscr{O}(k)$ equipped with the additional weight $\varphi=\alpha \log |P|^{2}$. We may assume that the open set $U=\{|\sigma|<1\}$ is such that $\operatorname{nil}(\Delta) \cap \bar{U}=\emptyset$, otherwise it suffices to multiply $\sigma$ by a large constant. This implies that the polynomial $Q=\prod P_{j}^{\left[k_{j} \alpha\right]}$ satisfies

$$
\int_{U}|Q|^{2} e^{-\varphi} d V=\int_{U} \prod\left|P_{j}\right|^{-2\left(k_{j} \alpha-\left[k_{j} \alpha\right]\right)} d V<+\infty
$$

Set $\omega=i c(\mathbb{O}(1))$. We have $\mathrm{id} d^{\prime} d^{\prime \prime} \log |P|^{2} \geqslant-i c(\mathbb{O}(D))=-D \omega$ by the Lelong-Poincaré equation, thus $\mathrm{i} \Theta\left(L_{\varphi}\right) \geqslant(k-\alpha D) \omega$. The desired curvature inequality $\mathrm{i} \Theta\left(L_{\varphi}\right) \geqslant(a+1+$ $\varepsilon) \mathrm{i} \Theta(E)$ is satisfied if $k-\alpha D \geqslant(a+1+\varepsilon)$. We thus take

$$
k=[\alpha D]+a+2 .
$$

The section $f \in H^{0}\left(U, K_{\mathbb{P}^{n}} \otimes L\right)=H^{0}(U, \mathscr{O}(k-n-1))$ is taken to be a multiple of $Q$ by some polynomial. This is possible provided that

$$
k-n-1 \geqslant \operatorname{deg} Q \quad \Longleftrightarrow \quad \alpha D+a+2-n-1 \geqslant \sum\left[k_{j} \alpha\right] \operatorname{deg} P_{j}
$$

or equivalently, as $D=\sum k_{j} \operatorname{deg} P_{j}$,

$$
\begin{equation*}
\sum\left(k_{j} \alpha-\left[k_{j} \alpha\right]\right) \operatorname{deg} P_{j} \geqslant n-a-1 \tag{8.10}
\end{equation*}
$$

Then we get $f \in H^{0}\left(U, K_{\mathbb{P}^{n}} \otimes L\right)$ such that $\int_{U}|f|^{2} e^{-\varphi} d V<+\infty$. Theorem 7.1 implies the existence of $F \in H^{0}\left(\mathbb{P}^{n}, K_{\mathbb{P}^{n}} \otimes L\right)$, i.e. of a polynomial $F$ of degree $k-n-1$, such that

$$
\int_{\mathbb{P}^{n}}|F|^{2} e^{-\varphi} d V=\int_{\mathbb{P}^{n}} \frac{|F|^{2}}{|P|^{2 \alpha}} d V<+\infty
$$

observe that $|\sigma|$ is bounded, for we are on a compact manifold. Near any $w \in S$, we have $|P(z)| \leqslant C|z-w|^{t_{2}}$, thus $|P(z)|^{2 \alpha} \leqslant C|z-w|^{2\left(t_{1}+n-1\right)}$. This implies that the above integral can converge only if $F$ vanishes at order $\geqslant t_{1}$ at each point $w \in S$. Therefore

$$
\omega_{t_{1}}(S) \leqslant \operatorname{deg} F=k-n-1=[\alpha D]+a+1-n \leqslant \alpha \omega_{t_{2}}(S)+a+1-n
$$

which is the desired inequality.
However, the above proof only works under the additional assumption (8.10). Assume on the contrary that

$$
\beta=\sum\left(k_{j} \alpha-\left[k_{j} \alpha\right]\right) \operatorname{deg} P_{j}<n-a-1 .
$$

Then the polynomial $Q$ has degree

$$
\sum\left[k_{j} \alpha\right] \operatorname{deg} P_{j}=\alpha \operatorname{deg} P-\beta=\alpha D-\beta
$$

and $Q$ vanishes at every point $w \in S$ with order

$$
\begin{aligned}
\operatorname{ord}_{w} Q & \geqslant \sum\left[k_{j} \alpha\right] \operatorname{ord}_{w} P_{j}=\alpha \sum k_{j} \operatorname{ord}_{w} P_{j}-\sum\left(k_{j} \alpha-\left[k_{j} \alpha\right]\right) \operatorname{ord}_{w} P_{j} \\
& \geqslant \alpha \operatorname{ord}_{w} P-\beta \geqslant \alpha t_{2}-\beta=t_{1}-(\beta-n+1)
\end{aligned}
$$

This implies $\operatorname{ord}_{w} Q \geqslant t_{1}-[\beta-n+1]$. As $[\beta-n+1]<n-a-1-n+1=-a \leqslant 0$, we can take a derivative of order $-[\beta-n+1]$ of $Q$ to get a polynomial $F$ with

$$
\operatorname{deg} F=\alpha D-\beta+[\beta-n+1] \leqslant \alpha D-n+1
$$

which vanishes at order $t_{1}$ on $S$. In this case, we obtain therefore

$$
\omega_{t_{1}}(S) \leqslant \alpha D-n+1=\frac{t_{1}+n-1}{t_{2}} \omega_{t_{2}}(S)-n+1
$$

and the proof of Th. 8.6 is complete.

## § 9. Skoda's $L^{\mathbf{2}}$ Estimates for Surjective Bundle Morphisms

Let $(X, \omega)$ be a Kähler manifold, $\operatorname{dim} X=n$, and $g: E \longrightarrow Q$ a holomorphic morphism of hermitian vector bundles over $X$. Assume in the first instance that $g$ is surjective. We are interested in conditions insuring for example that the induced morphism $g: H^{k}\left(X, K_{X} \otimes E\right) \longrightarrow H^{k}\left(X, K_{X} \otimes Q\right)$ is also surjective. For that purpose, it is natural to consider the subbundle $S=\operatorname{Ker} g \subset E$ and the exact sequence

$$
\begin{equation*}
0 \longrightarrow S \longrightarrow E \xrightarrow{g} Q \longrightarrow 0 \tag{9.1}
\end{equation*}
$$

Assume for the moment that $S$ and $Q$ are endowed with the metrics induced by that of $E$. Let $L$ be a line bundle over $X$. We consider the tensor product of sequence (9.1) by $L$ :

$$
\begin{equation*}
0 \longrightarrow S \otimes L \longrightarrow E \otimes L \xrightarrow{g} Q \otimes L \longrightarrow 0 . \tag{9.2}
\end{equation*}
$$

(9.3) Theorem. Let $k$ be an integer such that $0 \leqslant k \leqslant n$. Set $r=\operatorname{rk} E, q=\operatorname{rkQ}$, $s=\operatorname{rk} S=r-q$ and

$$
m=\min \{n-k, s\}=\min \{n-k, r-q\} .
$$

Assume that $(X, \omega)$ possesses also a complete Kähler metric $\widehat{\omega}$, that $E \geqslant_{m} 0$, and that $L \longrightarrow X$ is a hermitian line bundle such that

$$
\mathrm{i} \Theta(L)-(m+\varepsilon) \mathrm{i} \Theta(\operatorname{det} Q) \geqslant 0
$$

for some $\varepsilon>0$. Then for every $D^{\prime \prime}$-closed form $f$ of type $(n, k)$ with values in $Q \otimes L$ such that $\|f\|<+\infty$, there exists a $D^{\prime \prime}$-closed form $h$ of type $(n, k)$ with values in $E \otimes L$ such that $f=g \cdot h$ and

$$
\|h\|^{2} \leqslant(1+m / \varepsilon)\|f\|^{2} .
$$

The idea of the proof is essentially due to [Skoda 1978], who actually proved the special case $k=0$. The general case appeared in [Demailly 1982c].

Proof. Let $j: S \rightarrow E$ be the inclusion morphism, $g^{\star}: Q \rightarrow E$ and $j^{\star}: E \rightarrow S$ the adjoints of $g, j$, and

$$
D_{E}=\left(\begin{array}{cc}
D_{S} & -\beta^{\star} \\
\beta & D_{Q}
\end{array}\right), \quad \beta \in \mathscr{C}_{1,0}^{\infty}(X, \operatorname{hom}(S, Q)), \quad \beta^{\star} \in \mathscr{C}_{0,1}^{\infty}(X, \operatorname{hom}(Q, S))
$$

the matrix of $D_{E}$ with respect to the orthogonal splitting $E \simeq S \oplus Q$ (cf. §V-14). Then $g^{\star} f$ is a lifting of $f$ in $E \otimes L$. We shall try to find $h$ under the form

$$
h=g^{\star} f+j u, \quad u \in L_{n, k}^{2}(X, S \otimes L) .
$$

As the images of $S$ and $Q$ in $E$ are orthogonal, we have $|h|^{2}=|f|^{2}+|u|^{2}$ at every point of $X$. On the other hand $D_{Q \otimes L}^{\prime \prime} f=0$ by hypothesis and $D^{\prime \prime} g^{\star}=-j \circ \beta^{\star}$ by V-14.3 d), hence

$$
D_{E \otimes L}^{\prime \prime} h=-j\left(\beta^{\star} \wedge f\right)+j D_{S \otimes L}^{\prime \prime}=j\left(D_{S \otimes L}^{\prime \prime}-\beta^{\star} \wedge f\right)
$$

We are thus led to solve the equation

$$
\begin{equation*}
D_{S \otimes L}^{\prime \prime} u=\beta^{\star} \wedge f \tag{9.4}
\end{equation*}
$$

and for that, we apply Th. 4.5 to the $(n, k+1)$-form $\beta^{\star} \wedge f$. One observes now that the curvature of $S \otimes L$ can be expressed in terms of $\beta$. This remark will be used to prove:
(9.5) Lemma. $\left\langle A_{k}^{-1}\left(\beta^{\star} \wedge f\right),\left(\beta^{\star} \wedge f\right)\right\rangle \leqslant(m / \varepsilon)|f|^{2}$.

If the Lemma is taken for granted, Th. 4.5 yields a solution $u$ of (9.4) in $L_{n, q}^{2}(X, S \otimes L)$ such that $\|u\|^{2} \leqslant(m / \varepsilon)\|f\|^{2}$. As $\|h\|^{2}=\|f\|^{2}+\|u\|^{2}$, the proof of Th. 9.3 is complete.

Proof of Lemma 9.5.. Exactly as in the proof of Th. VII-10.3, formulas (V-14.6) and (V-14.7) yield

$$
\mathrm{i} \Theta(S) \geqslant{ }_{m} \mathrm{i} \beta^{\star} \wedge \beta, \quad \mathrm{i} \Theta(\operatorname{det} Q) \geqslant \operatorname{Tr}_{Q}\left(\mathrm{i} \beta \wedge \beta^{\star}\right)=\operatorname{Tr}_{S}\left(-\mathrm{i} \beta^{\star} \wedge \beta\right)
$$

Since $\mathscr{C}_{1,1}^{\infty}(X, \operatorname{Herm} S) \ni \Theta:=-\mathrm{i} \beta^{\star} \wedge \beta \geqslant_{\text {Grif }} 0$, Prop. VII-10.1 implies

$$
m \operatorname{Tr}_{S}\left(-\mathrm{i} \beta^{\star} \wedge \beta\right) \otimes \operatorname{Id}_{S}+\mathrm{i} \beta^{\star} \wedge \beta \geqslant_{m} 0
$$

From the hypothesis on the curvature of $L$ we get

$$
\begin{aligned}
\mathrm{i} \Theta(S \otimes L) & \geqslant_{m} \mathrm{i} \Theta(S) \otimes \operatorname{Id}_{L}+(m+\varepsilon) \mathrm{i} \Theta(\operatorname{det} Q) \otimes \operatorname{Id}_{S \otimes L} \\
& \geqslant_{m}\left(\mathrm{i} \beta^{\star} \wedge \beta+(m+\varepsilon) \operatorname{Tr}_{S}\left(-\mathrm{i} \beta^{\star} \wedge \beta\right) \otimes \operatorname{Id}_{S}\right) \otimes \operatorname{Id}_{L} \\
& \geqslant_{m}(\varepsilon / m)\left(-\mathrm{i} \beta^{\star} \wedge \beta\right) \otimes \operatorname{Id}_{S} \otimes \operatorname{Id}_{L}
\end{aligned}
$$

For any $v \in \Lambda^{n, k+1} T_{X}^{\star} \otimes S \otimes L$, Lemma VII-7.2 implies

$$
\begin{equation*}
\left\langle A_{k, S \otimes L} v, v\right\rangle \geqslant(\varepsilon / m)\left\langle-\mathrm{i} \beta^{\star} \wedge \beta \wedge \Lambda v, v\right\rangle, \tag{9.6}
\end{equation*}
$$

because $\operatorname{rk}(S \otimes L)=s$ and $m=\min \{n-k, s\}$. Let $\left(d z_{1}, \ldots, d z_{n}\right)$ be an orthonormal basis of $T_{X}^{\star}$ at a given point $x_{0} \in X$ and set

$$
\beta=\sum_{1 \leqslant j \leqslant n} d z_{j} \otimes \beta_{j}, \quad \beta_{j} \in \operatorname{hom}(S, Q)
$$

The adjoint of the operator $\beta^{\star} \wedge \bullet=\sum d \bar{z}_{j} \wedge \beta_{j}^{\star} \bullet$ is the contraction $\left.\beta\right\lrcorner \bullet$ defined by

$$
\left.\beta\lrcorner v=\sum \frac{\partial}{\partial \bar{z}_{j}}\right\lrcorner\left(\beta_{j} v\right)=\sum-\mathrm{i} d z_{j} \wedge \Lambda\left(\beta_{j} v\right)=-\mathrm{i} \beta \wedge \Lambda v .
$$

We get consequently $\left.\left\langle-\mathrm{i} \beta^{\star} \wedge \beta \wedge \Lambda v, v\right\rangle=\mid \beta\right\lrcorner\left. v\right|^{2}$ and (9.6) implies

$$
\left.\left|\left\langle\beta^{\star} \wedge f, v\right\rangle\right|^{2}=\left.|\langle f, \beta\lrcorner v\rangle\right|^{2} \leqslant|f|^{2} \mid \beta\right\lrcorner\left. v\right|^{2} \leqslant(m / \varepsilon)\left\langle A_{k, S \otimes L} v, v\right\rangle|f|^{2} .
$$

If $X$ has a plurisubharmonic exhaustion function $\psi$, we can select a convex increasing function $\chi \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$ and multiply the metric of $L$ by the weight $\exp (-\chi \circ \psi)$ in order to make the $L^{2}$ norm of $f$ converge. Theorem 9.3 implies therefore:
(9.7) Corollary. Let $(X, \omega)$ be a weakly pseudoconvex Kähler manifold, let $g: E \rightarrow Q$ be a surjective bundle morphism with $r=\mathrm{rk} E, q=\mathrm{rk} Q$, let $m=\min \{n-k, r-q\}$ and let $L \rightarrow X$ be a hermitian line bundle. Suppose that $E \geqslant_{m} 0$ and

$$
\mathrm{i} \Theta(L)-(m+\varepsilon) \mathrm{i} \Theta(\operatorname{det} Q) \geqslant 0
$$

for some $\varepsilon>0$. Then $g$ induces a surjective map

$$
H^{k}\left(X, K_{X} \otimes E \otimes L\right) \longrightarrow H^{k}\left(X, K_{X} \otimes Q \otimes L\right)
$$

The most remarkable feature of this result is that it does not require any strict positivity assumption on the curvature (for instance $E$ can be a flat bundle). A careful examination of the proof shows that it amounts to verify that the image of the coboundary morphism

$$
-\beta^{\star} \wedge \bullet: H^{k}\left(X, K_{X} \otimes Q \otimes L\right) \longrightarrow H^{k+1}\left(X, K_{X} \otimes S \otimes L\right)
$$

vanishes; however the cohomology group $H^{k+1}\left(X, K_{X} \otimes S \otimes L\right)$ itself does not vanish in general as it would do under a strict positivity assumption (cf. Th. VII-9.4).

We want now to get also estimates when $Q$ is endowed with a metric given a priori, that can be distinct from the quotient metric of $E$ by $g$. Then the map $g^{\star}\left(g g^{\star}\right)^{-1}: Q \rightarrow E$ is the lifting of $Q$ orthogonal to $S=\operatorname{Ker} g$. The quotient metric $|\bullet|^{\prime}$ on $Q$ is therefore defined in terms of the original metric $|\bullet|$ by

$$
|v|^{\prime 2}=\left|g^{\star}\left(g g^{\star}\right)^{-1} v\right|^{2}=\left\langle\left(g g^{\star}\right)^{-1} v, v\right\rangle=\operatorname{det}\left(g g^{\star}\right)^{-1}\left\langle\widetilde{g g^{\star}} v, v\right\rangle
$$

where $\widetilde{g g^{\star}} \in \operatorname{End}(Q)$ denotes the endomorphism of $Q$ whose matrix is the transposed of the comatrix of $g g^{\star}$. For every $w \in \operatorname{det} Q$, we find

$$
|w|^{\prime 2}=\operatorname{det}\left(g g^{\star}\right)^{-1}|w|^{2} .
$$

If $Q^{\prime}$ denotes the bundle $Q$ with the quotient metric, we get

$$
\mathrm{i} \Theta\left(\operatorname{det} Q^{\prime}\right)=\mathrm{i} \Theta(\operatorname{det} Q)+\mathrm{i} d^{\prime} d^{\prime \prime} \log \operatorname{det}\left(g g^{\star}\right) .
$$

In order that the hypotheses of Th. 9.3 be satisfied, we are led to define a new metric $\mid$ - $\left.\right|^{\prime}$ on $L$ by $|u|^{\prime 2}=|u|^{2}\left(\operatorname{det}\left(g g^{\star}\right)\right)^{-m-\varepsilon}$. Then

$$
\mathrm{i} \Theta\left(L^{\prime}\right)=\mathrm{i} \Theta(L)+(m+\varepsilon) \mathrm{i} d^{\prime} d^{\prime \prime} \log \operatorname{det}\left(g g^{\star}\right) \geqslant(m+\varepsilon) \mathrm{i} \Theta\left(\operatorname{det} Q^{\prime}\right)
$$

Theorem 9.3 applied to $\left(E, Q^{\prime}, L^{\prime}\right)$ can now be reformulated:
(9.8) Theorem. Let $X$ be a complete Kähler manifold equipped with a Kähler metric $\omega$ on $X$, let $E \rightarrow Q$ be a surjective morphism of hermitian vector bundles and let $L \rightarrow X$ be a hermitian line bundle. Set $r=\operatorname{rk} E, q=\operatorname{rk} Q$ and $m=\min \{n-k, r-q\}$ and suppose $E \geqslant_{m} 0$,

$$
\mathrm{i} \Theta(L)-(m+\varepsilon) \mathrm{i} \Theta(\operatorname{det} Q) \geqslant 0
$$

for some $\varepsilon>0$. Then for every $D^{\prime \prime}$-closed form $f$ of type $(n, k)$ with values in $Q \otimes L$ such that

$$
I=\int_{X}\left\langle\widetilde{g g^{\star}} f, f\right\rangle\left(\operatorname{det} g g^{\star}\right)^{-m-1-\varepsilon} d V<+\infty,
$$

there exists a $D^{\prime \prime}$-closed form $h$ of type $(n, k)$ with values in $E \otimes L$ such that $f=g \cdot h$ and

$$
\int_{X}|h|^{2}\left(\operatorname{det} g g^{\star}\right)^{-m-\varepsilon} d V \leqslant(1+m / \varepsilon) I .
$$

Our next goal is to extend Th. 9.8 in the case when $g: E \longrightarrow Q$ is only generically surjective; this means that the analytic set

$$
Y=\left\{x \in X ; g_{x}: E_{x} \longrightarrow Q_{x} \text { is not surjective }\right\}
$$

defined by the equation $\Lambda^{q} g=0$ is nowhere dense in $X$. Here $\Lambda^{q} g$ is a section of the bundle $\operatorname{hom}\left(\Lambda^{q} E, \operatorname{det} Q\right)$.
(9.9) Theorem. The existence statement and the estimates of Th. 9.8 remain true for a generically surjective morphism $g: E \rightarrow Q$ provided that $X$ is weakly pseudoconvex.

Proof. Apply Th. 9.8 to each relatively compact domain $X_{c} \backslash Y$ (these domains are complete Kähler by Lemma 7.2). From a sequence of solutions on $X_{c} \backslash Y$ we can extract a subsequence converging weakly on $X \backslash Y$ as c tends to $+\infty$. One gets a form $h$ satisfying the estimates, such that $D^{\prime \prime} h=0$ on $X \backslash Y$ and $f=g \cdot h$. In order to see that $D^{\prime \prime} h=0$ on $X$, it suffices to apply Lemma 7.3 and to observe that $h$ has $L_{\mathrm{loc}}^{2}$ coefficients on $X$ by our estimates.

A very special but interesting case is obtained for the trivial bundles $E=\Omega \times \mathbb{C}^{r}$, $Q=\Omega \times \mathbb{C}$ over a pseudoconvex open set $\Omega \subset \mathbb{C}^{n}$. Then the morphism $g$ is given by a $r$-tuple $\left(g_{1}, \ldots, g_{r}\right)$ of holomorphic functions on $\Omega$. Let us take $k=0$ and $L=\Omega \times \mathbb{C}$ with the metric given by a weight $e^{-\varphi}$. If we observe that $\widetilde{g g^{\star}}=\operatorname{Id}$ when $\operatorname{rk} Q=1$, Th. 9.8 applied on $X=\Omega \backslash g^{-1}(0)$ and Lemmas 7.2, 7.3 give:
(9.10) Theorem ([Skoda 1978]). Let $\Omega$ be a complete Kähler open subset of $\mathbb{C}^{n}$ and $\varphi$ a plurisubharmonic function on $\Omega$. Set $m=\min \{n, r-1\}$. Then for every holomorphic function $f$ on $\Omega$ such that

$$
I=\int_{\Omega \backslash Z}|f|^{2}|g|^{-2(m+1+\varepsilon)} e^{-\varphi} d V<+\infty
$$

where $Z=g^{-1}(0)$, there exist holomorphic functions $\left(h_{1}, \ldots, h_{r}\right)$ on $\Omega$ such that $f=$ $\sum g_{j} h_{j}$ and

$$
\int_{\Omega \backslash Y}|h|^{2}|g|^{-2(m+\varepsilon)} e^{-\varphi} d V \leqslant(1+m / \varepsilon) I .
$$

This last theorem can be used in order to obtain a quick solution of the Levi problem mentioned in §I-4. It can be used also to prove a result of [Diederich-Pflug 1981], relating the pseudoconvexity property and the existence of complete Kähler metrics for domains of $\mathbb{C}^{n}$.
(9.11) Theorem. Let $\Omega \subset \mathbb{C}^{n}$ be an open subset. Then:
a) $\Omega$ is a domain of holomorphy if and only if $\Omega$ is pseudoconvex ;
b) If $(\bar{\Omega})^{\circ}=\Omega$ and if $\Omega$ has a complete Kähler metric $\widehat{\omega}$, then $\Omega$ is pseudoconvex.

Note that statement b) can be false if the assumption $(\bar{\Omega})^{\circ}=\Omega$ is omitted: in fact $\mathbb{C}^{n} \backslash\{0\}$ is complete Kähler by Lemma 7.2 , but it is not pseudoconvex if $n \geqslant 2$.

Proof. b) By Th. I-4.12, it is enough to verify that $\Omega$ is a domain of holomorphy, i.e. that for every connected open subset $U$ such that $U \cap \partial \Omega \neq \emptyset$ and every connected component W of $U \cap \Omega$ there exists a holomorphic function $h$ on $\Omega$ such that $h_{\upharpoonright W}$ cannot be continued to $U$. Since $(\bar{\Omega})^{\circ}=\Omega$, the set $U \backslash \bar{\Omega}$ is not empty. We select $a \in U \backslash \bar{\Omega}$. Then the integral

$$
\int_{\Omega}|z-a|^{-2(n+\varepsilon)} d V(z)
$$

converges. By Th. 9.10 applied to $f(z)=1, g_{j}(z)=z_{j}-a_{j}$ and $\varphi=0$, there exist holomorphic functions $h_{j}$ on $\Omega$ such that $\sum\left(z_{j}-a_{j}\right) h_{j}(z)=1$. This shows that at least one of the functions $h_{j}$ cannot be analytically continued at $a \in U$.
a) Assume that $\Omega$ is pseudoconvex. Given any open connected set $U$ such that $U \cap \partial \Omega \neq \emptyset$, choose $a \in U \cap \partial \Omega$. By Th. I-4.14 c) the function

$$
\varphi(z)=(n+\varepsilon)\left(\log \left(1+|z|^{2}\right)-2 \log d(z, \mathrm{C} \Omega)\right)
$$

is plurisubharmonic on $\Omega$. Then the integral

$$
\int_{\Omega}|z-a|^{-2(n+\varepsilon)} e^{-\varphi(z)} d V(z) \leqslant \int_{\Omega}\left(1+|z|^{2}\right)^{-n-\varepsilon} d V(z)
$$

converges, and we conclude as for b).

## § 10. Application of Skoda's $L^{2}$ Estimates to Local Algebra

We apply here Th. 9.10 to the study of ideals in the ring $\mathscr{O}_{n}=\mathbb{C}\left\{z_{1}, \ldots, z_{n}\right\}$ of germs of holomorphic functions on $\left(\mathbb{C}^{n}, 0\right)$. Let $\mathscr{I}=\left(g_{1}, \ldots, g_{r}\right) \neq(0)$ be an ideal of $\mathfrak{C}_{n}$.
(10.1) Definition. Let $k \in \mathbb{R}_{+}$. We associate to $\mathscr{I}$ the following ideals:
a) the ideal $\overline{\mathscr{J}}^{(k)}$ of germs $u \in \mathscr{O}_{n}$ such that $|u| \leqslant C|g|^{k}$ for some constant $C \geqslant 0$, where $|g|^{2}=\left|g_{1}\right|^{2}+\cdots+\left|g_{r}\right|^{2}$.
b) the ideal $\widehat{\mathscr{I}}^{(k)}$ of germs $u \in \mathfrak{O}_{n}$ such that

$$
\int_{\Omega}|u|^{2}|g|^{-2(k+\varepsilon)} d V<+\infty
$$

on a small ball $\Omega$ centered at 0 , if $\varepsilon>0$ is small enough.
(10.2) Proposition. For all $k, l \in \mathbb{R}_{+}$we have
a) $\overline{\mathscr{I}}^{(k)} \subset \widehat{\mathscr{I}}^{(k)}$;
b) $\mathscr{I}^{k} \subset \overline{\mathscr{I}}^{(k)}$ if $k \in \mathbb{N}$;
c) $\overline{\mathscr{I}}^{(k)} \cdot \overline{\mathscr{I}}^{(l)} \subset \overline{\mathscr{I}}^{(k+l)}$;
d) $\overline{\mathscr{I}}^{(k)} . \widehat{\mathscr{I}}^{(l)} \subset \widehat{\mathscr{I}}^{(k+l)}$.

All properties are immediate from the definitions except a) which is a consequence of Lemma 8.2. Before stating the main result, we need a simple lemma.
(10.3) Lemma. If $\mathscr{I}=\left(g_{1}, \ldots, g_{r}\right)$ and $r>n$, we can find elements $\widetilde{g}_{1}, \ldots, \widetilde{g}_{n} \in \mathscr{I}$ such that $C^{-1}|g| \leqslant|\widetilde{g}| \leqslant C|g|$ on a neighborhood of 0 . Each $\widetilde{g}_{j}$ can be taken to be a linear combination

$$
\widetilde{g}_{j}=a_{j} . g=\sum_{1 \leqslant k \leqslant r} a_{j k} g_{k}, \quad a_{j} \in \mathbb{C}^{r} \backslash\{0\}
$$

where the coefficients $\left(\left[a_{1}\right], \ldots,\left[a_{n}\right]\right)$ are chosen in the complement of a proper analytic subset of $\left(\mathbb{P}^{r-1}\right)^{n}$.

It follows from the Lemma that the ideal $\mathscr{g}=\left(\widetilde{g}_{1}, \ldots, \widetilde{g}_{n}\right) \subset \mathscr{I}$ satisfies $\overline{\mathscr{g}}^{(k)}=\overline{\mathscr{I}}^{(k)}$ and $\widehat{\mathscr{g}}^{(k)}=\widehat{\mathscr{J}}^{(k)}$ for all $k$.

Proof. Assume that $g \in \mathscr{O}(\Omega)^{r}$. Consider the analytic subsets in $\Omega \times\left(\mathbb{P}^{r-1}\right)^{n}$ defined by

$$
\begin{aligned}
A & =\left\{\left(z,\left[w_{1}\right], \ldots,\left[w_{n}\right]\right) ; w_{j} . g(z)=0\right\} \\
A^{\star} & =\bigcup \text { irreducible components of } A \text { not contained in } g^{-1}(0) \times\left(\mathbb{P}^{r-1}\right)^{n} .
\end{aligned}
$$

For $z \notin g^{-1}(0)$ the fiber $A_{z}=\left\{\left(\left[w_{1}\right], \ldots,\left[w_{n}\right]\right) ; w_{j} . g(z)=0\right\}=A_{z}^{\star}$ is a product of $n$ hyperplanes in $\mathbb{P}^{r-1}$, hence $A \cap\left(\Omega \backslash g^{-1}(0)\right) \times\left(\mathbb{P}^{r-1}\right)^{n}$ is a fiber bundle with base $\Omega \backslash g^{-1}(0)$ and fiber $\left(\mathbb{P}^{r-2}\right)^{n}$. As $A^{\star}$ is the closure of this set in $\Omega \times\left(\mathbb{P}^{r-1}\right)^{n}$, we have

$$
\operatorname{dim} A^{\star}=n+n(r-2)=n(r-1)=\operatorname{dim}\left(\mathbb{P}^{r-1}\right)^{n}
$$

It follows that the zero fiber

$$
A_{0}^{\star}=A^{\star} \cap\left(\{0\} \times\left(\mathbb{P}^{r-1}\right)^{n}\right)
$$

is a proper subset of $\{0\} \times\left(\mathbb{P}^{r-1}\right)^{n}$. Choose $\left(a_{1}, \ldots, a_{n}\right) \in\left(\mathbb{C}^{r} \backslash\{0\}\right)^{n}$ such that $\left(0,\left[a_{1}\right], \ldots,\left[a_{n}\right]\right)$ is not in $A_{0}^{\star}$. By an easy compactness argument the set $A^{\star} \cap\left(\bar{B}(0, \varepsilon) \times\left(\mathbb{P}^{r-1}\right)^{n}\right)$ is disjoint from the neighborhood

$$
B(0, \varepsilon) \times \prod\left[B\left(a_{j}, \varepsilon\right)\right]
$$

of $\left(0,\left[a_{1}\right], \ldots,\left[a_{n}\right]\right)$ for $\varepsilon$ small enough. For $z \in B(0, \varepsilon)$ we have $\left|a_{j} . g(z)\right| \geqslant \varepsilon|g(z)|$ for some $j$, otherwise the inequality $\left|a_{j} . g(z)\right|<\varepsilon|g(z)|$ would imply the existence of $h_{j} \in \mathbb{C}^{r}$ with $\left|h_{j}\right|<\varepsilon$ and $a_{j} . g(z)=h_{j} . g(z)$. Since $g(z) \neq 0$, we would have

$$
\left(z,\left[a_{1}-h_{1}\right], \ldots,\left[a_{n}-h_{n}\right]\right) \in A^{\star} \cap\left(B(0, \varepsilon) \times\left(\mathbb{P}^{r-1}\right)^{n}\right)
$$

a contradiction. We obtain therefore

$$
\varepsilon|g(z)| \leqslant \max \left|a_{j} . g(z)\right| \leqslant\left(\max \left|a_{j}\right|\right)|g(z)| \quad \text { on } B(0, \varepsilon)
$$

(10.4) Theorem ([Briançon-Skoda 1974]). Set $p=\min \{n-1, r-1\}$. Then
a) $\widehat{\mathscr{I}}^{(k+1)}=\mathscr{I} \widehat{\mathscr{I}}^{(k)}=\overline{\mathscr{I}} \widehat{\mathscr{I}}^{(k)} \quad$ for $k \geqslant p$.
b) $\overline{\mathscr{I}}^{(k+p)} \subset \widehat{\mathscr{I}}^{(k+p)} \subset \mathscr{I}^{k} \quad$ for all $k \in \mathbb{N}$.

Proof. a) The inclusions $\mathscr{I} \widehat{\mathscr{I}}^{(k)} \subset \overline{\mathscr{I}} \widehat{\mathscr{I}}^{(k)} \subset \widehat{\mathscr{I}}^{(k+1)}$ are obvious thanks to Prop. 10.2, so we only have to prove that $\widehat{\mathscr{I}}^{(k+1)} \subset \mathscr{I}^{(k)}$. Assume first that $r \leqslant n$. Let $f \in \widehat{\mathscr{I}}^{(k+1)}$ be such that

$$
\int_{\Omega}|f|^{2}|g|^{-2(k+1+\varepsilon)} d V<+\infty
$$

For $k \geqslant p-1$, we can apply Th. 9.10 with $m=r-1$ and with the weight $\varphi=$ $(k-m) \log |g|^{2}$. Hence $f$ can be written $f=\sum g_{j} h_{j}$ with

$$
\int_{\Omega}|h|^{2}|g|^{-2(k+\varepsilon)} d V<+\infty
$$

thus $h_{j} \in \widehat{\mathscr{I}}^{(k)}$ and $f \in \mathscr{I} \widehat{\mathscr{I}}^{(k)}$. When $r>n$, Lemma 10.3 shows that there is an ideal $\mathscr{f} \subset \mathscr{I}$ with $n$ generators such that $\widehat{\mathscr{g}}^{(k)}=\widehat{\mathscr{I}}^{(k)}$. We find

$$
\widehat{\mathscr{I}}^{(k+1)}=\widehat{\mathscr{g}}^{(k+1)} \subset \mathscr{\mathscr { L }} \widehat{\mathscr{G}}^{(k)} \subset \mathscr{I} \widehat{\mathscr{I}}^{(k)} \quad \text { for } k \geqslant n-1 .
$$

b) Property a) implies inductively $\widehat{\mathscr{I}}^{(k+p)}=\mathscr{I}^{k} \widehat{\mathscr{I}}^{(p)}$ for all $k \in \mathbb{N}$. This gives in particular $\widehat{\mathscr{I}}^{(k+p)} \subset \mathscr{I}^{k}$.

## (10.5) Corollary.

a) The ideal $\overline{\mathscr{I}}$ is the integral closure of $\mathscr{I}$, i.e. by definition the set of germs $u \in \mathscr{O}_{n}$ which satisfy an equation

$$
u^{d}+a_{1} u^{d-1}+\cdots+a_{d}=0, \quad a_{s} \in \mathscr{I}^{s}, \quad 1 \leqslant s \leqslant d
$$

b) Similarly, $\overline{\mathscr{I}}^{(k)}$ is the set of germs $u \in \mathscr{O}_{n}$ which satisfy an equation

$$
u^{d}+a_{1} u^{d-1}+\cdots+a_{d}=0, \quad a_{s} \in I^{k s s}, \quad 1 \leqslant s \leqslant d
$$

where $] t$ [ denotes the smallest integer $\geqslant t$.

As the ideal $\overline{\mathscr{I}}^{(k)}$ is finitely generated, property b) shows that there always exists a rational number $l \geqslant k$ such that $\overline{\mathscr{I}}^{(l)}=\overline{\mathscr{I}}^{(k)}$.

Proof. a) If $u \in \mathscr{O}_{n}$ satisfies a polynomial equation with coefficients $a_{s} \in \mathscr{J}^{s}$, then clearly $\left|a_{s}\right| \leqslant C_{s}|g|^{s}$ and Lemma II-4.10 implies $|u| \leqslant C|g|$.

Conversely, assume that $u \in \overline{\mathscr{I}}$. The ring $\widehat{O}_{n}$ is Noetherian, so the ideal $\widehat{\mathscr{G}}(p)$ has a finite number of generators $v_{1}, \ldots, v_{N}$. For every $j$ we have $u v_{j} \in \overline{\mathscr{I}} \widehat{\mathscr{I}}^{(p)}=\mathscr{I} \widehat{\mathscr{I}}^{(p)}$, hence there exist elements $b_{j k} \in \mathscr{I}$ such that

$$
u v_{j}=\sum_{1 \leqslant k \leqslant N} b_{j k} v_{k}
$$

The matrix $\left(u \delta_{j k}-b_{j k}\right)$ has the non zero vector $\left(v_{j}\right)$ in its kernel, thus $u$ satisfies the equation $\operatorname{det}\left(u \delta_{j k}-b_{j k}\right)=0$, which is of the required type.
b) Observe that $v_{1}, \ldots, v_{N}$ satisfy simultaneously some integrability condition $\int_{\Omega}\left|v_{j}\right|^{-2(p+\varepsilon)}<+\infty$, thus $\widehat{\mathscr{I}}^{(p)}=\widehat{\mathscr{I}}^{(p+\eta)}$ for $\eta \in\left[0, \varepsilon\left[\right.\right.$. Let $u \in \overline{\mathscr{I}}^{(k)}$. For every integer $m \in \mathbb{N}$ we have

$$
u^{m} v_{j} \in \overline{\mathscr{I}}^{(k m)} \widehat{\mathscr{I}}^{(p+\eta)} \subset \widehat{\mathscr{I}}^{(k m+\eta+p)} .
$$

If $k \notin \mathbb{Q}$, we can find $m$ such that $d(k m+\varepsilon / 2, \mathbb{Z})<\varepsilon / 2$, thus $k m+\eta \in \mathbb{N}$ for some $\eta \in] 0, \varepsilon[$. If $k \in \mathbb{Q}$, we take $m$ such that $k m \in \mathbb{N}$ and $\eta=0$. Then

$$
u^{m} v_{j} \in \widehat{\mathscr{I}}^{(N+p)}=\mathscr{I}^{N} \widehat{\mathscr{I}}^{(p)} \quad \text { with } \quad N=k m+\eta \in \mathbb{N},
$$

and the reasoning made in a) gives $\operatorname{det}\left(u^{m} \delta_{j k}-b_{j k}\right)=0$ for some $b_{j k} \in \mathscr{I}^{N}$. This is an equation of the type described in b ), where the coefficients $a_{s}$ vanish when $s$ is not a multiple of $m$ and $\left.a_{m s} \in \mathscr{I}^{N s} \subset \mathscr{I}\right] k m s[$.

Let us mention that Briançon and Skoda's result 10.4 b ) is optimal for $k=1$. Take for example $\mathscr{I}=\left(g_{1}, \ldots, g_{r}\right)$ with $g_{j}(z)=z_{j}^{r}, 1 \leqslant j \leqslant r$, and $f(z)=z_{1} \ldots z_{r}$. Then $|f| \leqslant C|g|$ and 10.4 b ) yields $f^{r} \in \mathscr{I}$; however, it is easy to verify that $f^{r-1} \notin \mathscr{I}$. The theorem also gives an answer to the following conjecture made by J. Mather.
(10.6) Corollary. Let $f \in \mathscr{O}_{n}$ and $\mathscr{I}_{f}=\left(z_{1} \partial f / \partial z_{1}, \ldots, z_{n} \partial f / \partial z_{n}\right)$. Then $f \in \overline{\mathscr{I}}_{f}$, and for every integer $k \geqslant 0, f^{k+n-1} \in \mathscr{J}_{f}^{k}$.

The Corollary is also optimal for $k=1$ : for example, one can verify that the function $f(z)=\left(z_{1} \ldots z_{n}\right)^{3}+z_{1}^{3 n-1}+\ldots+z_{n}^{3 n-1}$ is such that $f^{n-1} \notin \mathscr{I}_{f}$.

Proof. Set $g_{j}(z)=z_{j} \partial f / \partial z_{j}, 1 \leqslant j \leqslant n$. By 10.4 b$)$, it suffices to show that $|f| \leqslant C|g|$. For every germ of analytic curve $\mathbb{C} \ni t \longmapsto \gamma(t), \gamma \not \equiv 0$, the vanishing order of $f \circ \gamma(t)$ at $t=0$ is the same as that of

$$
t \frac{d(f \circ \gamma)}{d t}=\sum_{1 \leqslant j \leqslant n} t \gamma_{j}^{\prime}(t) \frac{\partial f}{\partial z_{j}}(\gamma(t))
$$

We thus obtain

$$
|f \circ \gamma(t)| \leqslant C_{1}|t|\left|\frac{d(f \circ \gamma)}{d t}\right| \leqslant C_{2} \sum_{1 \leqslant j \leqslant n}\left|t \gamma_{j}^{\prime}(t)\right|\left|\frac{\partial f}{\partial z_{j}}(\gamma(t))\right| \leqslant C_{3}|g \circ \gamma(t)|
$$

and conclude by the following elementary lemma.
(10.7) Lemma. Let $f, g_{1}, \ldots, g_{r} \in \mathbb{O}_{n}$ be germs of holomorphic functions vanishing at 0 . Then we have $|f| \leqslant C|g|$ for some constant $C$ if and only if for every germ of analytic curve $\gamma$ through 0 there exists a constant $C_{\gamma}$ such that $|f \circ \gamma| \leqslant C_{\gamma}|g \circ \gamma|$.

Proof. If the inequality $|f| \leqslant C|g|$ does not hold on any neighborhood of 0 , the germ of analytic set $(A, 0) \subset\left(\mathbb{C}^{n+r}, 0\right)$ defined by

$$
g_{j}(z)=f(z) z_{n+j}, \quad 1 \leqslant j \leqslant r,
$$

contains a sequence of points $\left(z_{\nu}, g_{j}\left(z_{\nu}\right) / f\left(z_{\nu}\right)\right)$ converging to 0 as $\nu$ tends to $+\infty$, with $f\left(z_{\nu}\right) \neq 0$. Hence $(A, 0)$ contains an irreducible component on which $f \not \equiv 0$ and there is a germ of curve $\widetilde{\gamma}=\left(\gamma, \gamma_{n+j}\right):(\mathbb{C}, 0) \rightarrow\left(\mathbb{C}^{n+r}, 0\right)$ contained in $(A, 0)$ such that $f \circ \gamma \not \equiv 0$. We get $g_{j} \circ \gamma=(f \circ \gamma) \gamma_{n+j}$, hence $|g \circ \gamma(t)| \leqslant C|t||f \circ \gamma(t)|$ and the inequality $|f \circ \gamma| \leqslant C_{\gamma}|g \circ \gamma|$ does not hold.

## $\S$ 11. Integrability of Almost Complex Structures

Let $M$ be a $\mathscr{C}^{\infty}$ manifold of real dimension $m=2 n$. An almost complex structure on $M$ is by definition an endomorphism $J \in \operatorname{End}(T M)$ of class $\mathscr{C}^{\infty}$ such that $J^{2}=-\mathrm{Id}$. Then $T M$ becomes a complex vector bundle for which the scalar multiplication by $i$ is given by $J$. The pair $(M, J)$ is said to be an almost complex manifold. For such a manifold, the complexified tangent space $T_{\mathbb{C}} M=\mathbb{C} \otimes_{\mathbb{R}} T M$ splits into conjugate complex subspaces

$$
\begin{equation*}
T_{\mathbb{C}} M=T^{1,0} M \oplus T^{0,1} M, \quad \operatorname{dim}_{\mathbb{C}} T^{1,0} M=\operatorname{dim}_{\mathbb{C}} T^{0,1} M=n \tag{11.1}
\end{equation*}
$$

where $T^{1,0} M, T^{0,1} M \subset T_{\mathbb{C}} M$ are the eigenspaces of $\operatorname{Id} \otimes J$ corresponding to the eigenvalues $i$ and $-i$. The complexified exterior algebra $\mathbb{C} \otimes_{\mathbb{R}} \Lambda^{\bullet} T^{\star} M=\Lambda^{\bullet} T_{\mathbb{C}}^{\star} M$ has a corresponding splitting

$$
\begin{equation*}
\Lambda^{k} T_{\mathbb{C}}^{\star} M=\bigoplus_{p+q=k} \Lambda^{p, q} T_{\mathbb{C}}^{\star} M \tag{11.2}
\end{equation*}
$$

where we denote by definition

$$
\begin{equation*}
\Lambda^{p, q} T_{\mathbb{C}}^{\star} M=\Lambda^{p}\left(T^{1,0} M\right)^{\star} \otimes_{\mathbb{C}} \Lambda^{q}\left(T^{0,1} M\right)^{\star} \tag{11.3}
\end{equation*}
$$

As for complex manifolds, we let $C_{p, q}^{s}(M, E)$ be the space of differential forms of class $C^{s}$ and bidegree $(p, q)$ on $M$ with values in a complex vector bundle $E$. There is a natural antisymmetric bilinear map

$$
\theta: \mathscr{C}^{\infty}\left(M, T^{1,0} M\right) \times \mathscr{C}^{\infty}\left(M, T^{1,0} M\right) \longrightarrow \mathscr{C}^{\infty}\left(M, T^{0,1} M\right)
$$

which associates to a pair $(\xi, \eta)$ of $(1,0)$-vector fields the $(0,1)$-component of the Lie bracket $[\xi, \eta]$. Since

$$
[\xi, f \eta]=f[\xi, \eta]+(\xi . f) \eta, \quad \forall f \in \mathscr{C}^{\infty}(M, \mathbb{C})
$$

we see that $\theta(\xi, f \eta)=f \theta(\xi, \eta)$. It follows that $\theta$ is in fact a $(2,0)$-form on $M$ with values in $T^{0,1} M$.

If $M$ is a complex analytic manifold and $J$ its natural almost complex structure, we have in fact $\theta=0$, because $\left[\partial / \partial z_{j}, \partial / \partial z_{k}\right]=0,1 \leqslant j, k \leqslant n$, for any holomorphic local coordinate system $\left(z_{1}, \ldots, z_{n}\right)$.
(11.4) Definition. The form $\theta \in \mathscr{C}_{2,0}^{\infty}\left(M, T^{0,1} M\right)$ is called the torsion form of $J$. The almost complex structure $J$ is said to be integrable if $\theta=0$.
(11.5) Example. If $M$ is of real dimension $m=2$, every almost complex structure is integrable, because $n=1$ and alternate ( 2,0 )-forms must be zero. Assume that $M$ is a smooth oriented surface. To any Riemannian metric $g$ we can associate the endomorphism $J \in \operatorname{End}(T M)$ equal to the rotation of $+\pi / 2$. A change of orientation changes $J$ into the conjugate structure $-J$. Conversely, if $J$ is given, $T M$ is a complex line bundle, so $M$ is oriented, and a Riemannian metric $g$ is associated to $J$ if and only if $g$ is $J$-hermitian. As a consequence, there is a one-to-one correspondence between conformal classes of Riemannian metrics on $M$ and almost complex structures corresponding to a given orientation.

If $(M, J)$ is an almost complex manifold and $u \in \mathscr{C}_{p, q}^{\infty}(M, \mathbb{C})$, we let $d^{\prime} u, d^{\prime \prime} u$ be the components of type $(p+1, q)$ and $(p, q+1)$ in the exterior derivative $d u$. Let $\left(\xi_{1}, \ldots, \xi_{n}\right)$ be a frame of $T^{1,0} M_{\upharpoonright \Omega}$. The torsion form $\theta$ can be written

$$
\theta=\sum_{1 \leqslant j \leqslant n} \alpha_{j} \otimes \bar{\xi}_{j}, \quad \alpha_{j} \in \mathscr{C}_{2,0}^{\infty}(\Omega, \mathbb{C}) .
$$

Then $\theta$ yields conjugate operators $\theta^{\prime}, \theta^{\prime \prime}$ on $\Lambda^{\bullet} T_{\mathbb{C}}^{\star} M$ such that

$$
\begin{equation*}
\left.\left.\theta^{\prime} u=\sum_{1 \leqslant j \leqslant n} \alpha_{j} \wedge\left(\bar{\xi}_{j}\right\lrcorner u\right), \quad \theta^{\prime \prime} u=\sum_{1 \leqslant j \leqslant n} \bar{\alpha}_{j} \wedge\left(\xi_{j}\right\lrcorner u\right) . \tag{11.6}
\end{equation*}
$$

If $u$ is of bidegree $(p, q)$, then $\theta^{\prime} u$ and $\theta^{\prime \prime} u$ are of bidegree $(p+2, q-1)$ and $(p-1, q+2)$. It is clear that $\theta^{\prime}, \theta^{\prime \prime}$ are derivations, i.e.

$$
\theta^{\prime}(u \wedge v)=\left(\theta^{\prime} u\right) \wedge v+(-1)^{\operatorname{deg} u} u \wedge\left(\theta^{\prime} v\right)
$$

for all smooth forms $u, v$, and similarly for $\theta^{\prime \prime}$.
(11.7) Proposition. We have $d=d^{\prime}+d^{\prime \prime}-\theta^{\prime}-\theta^{\prime \prime}$.

Proof. Since all operators occuring in the formula are derivations, it is sufficient to check the formula for forms of degree 0 or 1 . If $u$ is of degree 0 , the result is obvious because $\theta^{\prime} u=\theta^{\prime \prime} u=0$ and $d u$ can only have components of types $(1,0)$ or $(0,1)$. If $u$ is a 1 -form and $\xi, \eta$ are complex vector fields, we have

$$
d u(\xi, \eta)=\xi \cdot u(\eta)-\eta \cdot d u(\xi)-u([\xi, \eta]) .
$$

When $u$ is of type $(0,1)$ and $\xi, \eta$ of type $(1,0)$, we find

$$
(d u)^{2,0}(\xi, \eta)=-u(\theta(\xi, \eta))
$$

thus $(d u)^{2,0}=-\theta^{\prime} u$, and of course $(d u)^{1,1}=d^{\prime} u,(d u)^{0,2}=d^{\prime \prime} u, \theta^{\prime \prime} u=0$ by definition. The case of a ( 1,0 )-form $u$ follows by conjugation.

Proposition 11.7 shows that $J$ is integrable if and only if $d=d^{\prime}+d^{\prime \prime}$. In this case, we infer immediately

$$
d^{\prime 2}=0, \quad d^{\prime} d^{\prime \prime}+d^{\prime \prime} d^{\prime}=0, \quad d^{\prime \prime 2}=0
$$

For an integrable almost complex structure, we thus have the same formalism as for a complex analytic structure, and indeed we shall prove:
(11.8) Newlander-Nirenberg theorem (1957). Every integrable almost complex structure $J$ on $M$ is defined by a unique analytic structure.

The proof we shall give follows rather closely that of [Hörmander 1966], which was itself based on previous ideas of [Kohn 1963, 1964]. A function $f \in C^{1}(\Omega, \mathbb{C}), \Omega \subset M$, is said to be $J$-holomorphic if $d^{\prime \prime} f=0$. Let $f_{1}, \ldots, f_{p} \in C^{1}(\Omega, \mathbb{C})$ and let $h$ be a function of class $C^{1}$ on an open subset of $\mathbb{C}^{p}$ containing the range of $f=\left(f_{1}, \ldots, f_{p}\right)$. An easy computation gives

$$
\begin{equation*}
d^{\prime \prime}(h \circ f)=\sum_{1 \leqslant j \leqslant p}\left(\frac{\partial h}{\partial z_{j}} \circ f\right) d^{\prime \prime} f_{j}+\left(\frac{\partial h}{\partial \bar{z}_{j}} \circ f\right) \overline{d^{\prime} f_{j}}, \tag{11.9}
\end{equation*}
$$

in particular $h \circ f$ is $J$-holomorphic as soon as $f_{1}, \ldots, f_{p}$ are $J$-holomorphic and $h$ holomorphic in the usual sense.

Constructing a complex analytic structure on $M$ amounts to show the existence of $J$ holomorphic complex coordinates $\left(z_{1}, \ldots, z_{n}\right)$ on a neighborhood $\Omega$ of every point $a \in M$. Formula (11.9) then shows that all coordinate changes $h:\left(z_{k}\right) \mapsto\left(w_{k}\right)$ are holomorphic in the usual sense, so that $M$ is furnished with a complex analytic atlas. The uniqueness of the analytic structure associated to $J$ is clear, since the holomorphic functions are characterized by the condition $d^{\prime \prime} f=0$. In order to show the existence, we need a lemma.
(11.10) Lemma. For every point $a \in M$ and every integer $s \geqslant 1$, there exist $\mathscr{C}^{\infty}$ complex coordinates $\left(z_{1}, \ldots, z_{n}\right)$ centered at a such that

$$
d^{\prime \prime} z_{j}=\mathrm{O}\left(|z|^{s}\right), \quad 1 \leqslant j \leqslant n .
$$

Proof. By induction on $s$. Let $\left(\xi_{1}^{\star}, \ldots, \xi_{n}^{\star}\right)$ be a basis of $\Lambda^{1,0} T_{\mathbb{C}}^{\star} M$. One can find complex functions $z_{j}$ such that $d z_{j}(a)=\xi_{j}^{\star}$, i.e.

$$
d^{\prime} z_{j}(a)=\xi_{j}^{\star}, \quad d^{\prime \prime} z_{j}(a)=0
$$

Then $\left(z_{1}, \ldots, z_{n}\right)$ satisfy the conclusions of the Lemma for $s=1$. If $\left(z_{1}, \ldots, z_{n}\right)$ are already constructed for the integer $s$, we have a Taylor expansion

$$
d^{\prime \prime} z_{j}=\sum_{1 \leqslant k \leqslant n} P_{j k}(z, \bar{z}) \overline{d^{\prime} z_{k}}+\mathrm{O}\left(|z|^{s+1}\right)
$$

where $P_{j k}(z, w)$ is a homogeneous polynomial in $(z, w) \in \mathbb{C}^{n} \times \mathbb{C}^{n}$ of total degree $s$. As $J$ is integrable, we have

$$
\begin{aligned}
0=d^{\prime \prime 2} z_{j} & =\sum_{1 \leqslant k, l \leqslant n} \frac{\partial P_{j k}}{\partial z_{l}} d^{\prime \prime} z_{l} \wedge \overline{d^{\prime} z_{k}}+\frac{\partial P_{j k}}{\partial \bar{z}_{l}} \overline{d^{\prime} z_{l}} \wedge \overline{d^{\prime} z_{k}}+\mathrm{O}\left(|z|^{s}\right) \\
& =\sum_{1 \leqslant k<l \leqslant n}\left[\frac{\partial P_{j k}}{\partial \bar{z}_{l}}-\frac{\partial P_{j l}}{\partial \bar{z}_{k}}\right] \overline{d^{\prime} z_{l}} \wedge \overline{d^{\prime} z_{k}}+\mathrm{O}\left(|z|^{s}\right)
\end{aligned}
$$

because $\partial P_{j k} / \partial z_{l}$ is of degree $s-1$ and $d^{\prime \prime} z_{l}=\mathrm{O}\left(|z|^{s}\right)$. Since the polynomial between brackets is of degree $s-1$, we must have

$$
\frac{\partial P_{j k}}{\partial \bar{z}_{l}}-\frac{\partial P_{j l}}{\partial \bar{z}_{k}}=0, \quad \forall j, k, l .
$$

We define polynomials $Q_{j}$ of degree $s+1$

$$
Q_{j}(z, \bar{z})=\int_{0}^{1} \sum_{1 \leqslant l \leqslant n} \bar{z}_{l} P_{j l}(z, t \bar{z}) d t
$$

Trivial computations show that

$$
\begin{aligned}
\frac{\partial Q_{j}}{\partial \bar{z}_{k}} & =\int_{0}^{1}\left(P_{j k}+\sum_{1 \leqslant l \leqslant n} \bar{z}_{l} \frac{\partial P_{j l}}{\partial \bar{z}_{k}}\right)(z, t \bar{z}) d t \\
& =\int_{0}^{1} \frac{d}{d t}\left[t P_{j k}(z, t \bar{z})\right] d t=P_{j k}(z, \bar{z}), \\
d^{\prime \prime}\left(z_{j}-Q_{j}(z, \bar{z})\right) & =d^{\prime \prime} z_{j}-\sum_{1 \leqslant k \leqslant n} \frac{\partial Q_{j}}{\partial \bar{z}_{k}} \overline{d^{\prime} z_{k}}-\sum_{1 \leqslant k \leqslant n} \frac{\partial Q_{j}}{\partial z_{k}} d^{\prime \prime} z_{k} \\
& =-\sum_{1 \leqslant k \leqslant n} \frac{\partial Q_{j}}{\partial z_{k}} d^{\prime \prime} z_{k}+\mathrm{O}\left(|z|^{s+1}\right)=\mathrm{O}\left(|z|^{s+1}\right)
\end{aligned}
$$

because $\partial Q_{j} / \partial z_{k}$ is of degree $s$ and $d^{\prime \prime} z_{l}=\mathrm{O}(|z|)$. The new coordinates

$$
\widetilde{z}_{j}=z_{j}-Q_{j}(z, \bar{z}), \quad 1 \leqslant j \leqslant n
$$

fulfill the Lemma at step $s+1$.
All usual notions defined on complex analytic manifolds can be extended to integrable almost complex manifolds. For example, a smooth function $\varphi$ is said to be strictly plurisubharmonic if $\mathrm{i} d^{\prime} d^{\prime \prime} \varphi$ is a positive definite (1,1)-form. Then $\omega=\mathrm{i} d^{\prime} d^{\prime \prime} \varphi$ is a Kähler metric on $(M, J)$.

In this context, all $L^{2}$ estimates proved in the previous paragraphs still apply to an integrable almost complex manifold; remember that the proof of the Bochner-KodairaNakano identity used only Taylor developments of order $\leqslant 2$, and the coordinates given by Lemma 11.10 work perfectly well for that purpose. In particular, Th. 6.5 is still valid.
(11.11) Lemma. Let $\left(z_{1}, \ldots, z_{n}\right)$ be coordinates centered at a point $a \in M$ with $d^{\prime \prime} z_{j}=$ $\mathrm{O}\left(|z|^{s}\right), s \geqslant 3$. Then the functions

$$
\left.\left.\psi(z)=|z|^{2}, \quad \varphi_{\varepsilon}(z)=|z|^{2}+\log \left(|z|^{2}+\varepsilon^{2}\right), \quad \varepsilon \in\right] 0,1\right]
$$

are strictly plurisubharmonic on a small ball $|z|<r_{0}$.
Proof. We have

$$
\mathrm{i} d^{\prime} d^{\prime \prime} \psi=\mathrm{i} \sum_{1 \leqslant j \leqslant n} d^{\prime} z_{j} \wedge \overline{d^{\prime} z_{j}}+d^{\prime} \bar{z}_{j} \wedge d^{\prime \prime} z_{j}+z_{j} d^{\prime} d^{\prime \prime} \bar{z}_{j}+\bar{z}_{j} d^{\prime} d^{\prime \prime} z_{j}
$$

The last three terms are $\mathrm{O}\left(|z|^{s}\right)$ and the first one is positive definite at $z=0$, so the result is clear for $\psi$. Moreover

$$
\begin{aligned}
\mathrm{i} d^{\prime} d^{\prime \prime} \varphi_{\varepsilon}=\mathrm{i} d^{\prime} d^{\prime \prime} \psi & +i \frac{\left(|z|^{2}+\varepsilon^{2}\right) \sum d^{\prime} z_{j} \wedge \overline{d^{\prime} z_{j}}-\sum \bar{z}_{j} d^{\prime} z_{j} \wedge \overline{\sum \bar{z}_{j} d^{\prime} z_{j}}}{\left(|z|^{2}+\varepsilon^{2}\right)^{2}} \\
& +\frac{\mathrm{O}\left(|z|^{s}\right)}{|z|^{2}+\varepsilon^{2}}+\frac{\mathrm{O}\left(|z|^{s+2}\right)}{\left(|z|^{2}+\varepsilon^{2}\right)^{2}}
\end{aligned}
$$

We observe that the first two terms are positive definite, whereas the remainder is $\mathrm{O}(|z|)$ uniformly in $\varepsilon$.

Proof of theorem 11.8.. With the notations of the previous lemmas, consider the pseudoconvex open set

$$
\Omega=\{|z|<r\}=\left\{\psi(z)-r^{2}<0\right\}, \quad r<r_{0},
$$

endowed with the Kähler metric $\omega=\mathrm{id} d^{\prime} d^{\prime \prime} \psi$. Let $h \in \mathscr{D}(\Omega)$ be a cut-off function with $0 \leqslant h \leqslant 1$ and $h=1$ on a neighborhood of $z=0$. We apply Th. 6.5 to the ( 0,1 )-forms

$$
g_{j}=d^{\prime \prime}\left(z_{j} h(z)\right) \in \mathscr{C}_{0,1}^{\infty}(\Omega, \mathbb{C})
$$

for the weight

$$
\varphi(z)=A|z|^{2}+(n+1) \log |z|^{2}=\lim _{\varepsilon \rightarrow 0} A|z|^{2}+(n+1) \log \left(|z|^{2}+\varepsilon^{2}\right)
$$

Lemma 11.11 shows that $\varphi$ is plurisubharmonic for $A \geqslant n+1$, and for $A$ large enough we obtain

$$
\mathrm{id} d^{\prime} d^{\prime \prime} \varphi+\operatorname{Ricci}(\omega) \geqslant \omega \quad \text { on } \Omega
$$

By Remark (6.8) we get a function $f_{j}$ such that $d^{\prime \prime} f_{j}=g_{j}$ and

$$
\int_{\Omega}\left|f_{j}\right|^{2} e^{-\varphi} d V \leqslant \int_{\Omega}\left|g_{j}\right|^{2} e^{-\varphi} d V
$$

As $g_{j}=d^{\prime \prime} z_{j}=\mathrm{O}\left(|z|^{s}\right)$ and $e^{-\varphi}=\mathrm{O}\left(|z|^{-2 n-2}\right)$ near $z=0$, the integral of $g_{j}$ converges provided that $s \geqslant 2$. Then $\int\left|f_{j}(z)\right|^{2}|z|^{-2 n-2} d V$ converges also at $z=0$. Since the solution $f_{j}$ is smooth, we must have $f_{j}(0)=d f_{j}(0)=0$. We set

$$
\widetilde{z}_{j}=z_{j} h(z)-f_{j}, \quad 1 \leqslant j \leqslant n .
$$

Then $\widetilde{z}_{j}$ is $J$-holomorphic and $d \widetilde{z}_{j}(0)=d z_{j}(0)$, so $\left(z_{1}, \ldots, z_{n}\right)$ is a $J$-holomorphic coordinate system at $z=0$.

In particular, any Riemannian metric on an oriented 2-dimensional real manifold defines a unique analytic structure. This fact will be used in order to obtain a simple proof of the well-known:
(11.12) Uniformization theorem. Every simply connected Riemann surface $X$ is biholomorphic either to $\mathbb{P}^{1}, \mathbb{C}$ or the unit disk $\Delta$.

Proof. We will merely use the fact that $H^{1}(X, \mathbb{R})=0$. If $X$ is compact, then $X$ is a complex curve of genus 0 , so $X \simeq \mathbb{P}^{1}$ by Th. VI-14.16. On the other hand, the elementary Riemann mapping theorem says that an open set $\Omega \subset \mathbb{C}$ with $H^{1}(\Omega, \mathbb{R})=0$ is either equal to $\mathbb{C}$ or biholomorphic to the unit disk. Thus, all we have to show is that a non compact Riemann surface $X$ with $H^{1}(X, \mathbb{R})=0$ can be embedded in the complex plane $\mathbb{C}$.

Let $\Omega_{\nu}$ be an exhausting sequence of relatively compact connected open sets with smooth boundary in $X$. We may assume that $X \backslash \Omega_{\nu}$ has no relatively compact connected components, otherwise we "fill the holes" of $\Omega_{\nu}$ by taking the union with all such components. We let $Y_{\underline{\nu}}$ be the double of the manifold with boundary $\left(\bar{\Omega}_{\nu}, \partial \Omega_{\nu}\right)$, i.e. the union of two copies of $\bar{\Omega}_{\nu}$ with opposite orientations and the boundaries identified. Then $Y_{\nu}$ is a compact oriented surface without boundary.
(11.13) Lemma. We have $H^{1}\left(Y_{\nu}, \mathbb{R}\right)=0$.

Proof. Let us first compute $H_{c}^{1}\left(\Omega_{\nu}, \mathbb{R}\right)$. Let $u$ be a closed 1-form with compact support in $\Omega_{\nu}$. By Poincaré duality $H_{c}^{1}(X, \mathbb{R})=0$, so $u=d f$ for some function $f \in \mathscr{D}(X)$. As $d f=0$ on a neighborhood of $X \backslash \Omega_{\nu}$ and as all connected components of this set are non compact, $f$ must be equal to the constant zero near $X \backslash \Omega_{\nu}$. Hence $u=d f$ is the zero class in $H_{c}^{1}\left(\Omega_{\nu}, \mathbb{R}\right)$ and we get $H_{c}^{1}\left(\Omega_{\nu}, \mathbb{R}\right)=H^{1}\left(\Omega_{\nu}, \mathbb{R}\right)=0$. The exact sequence of the pair $\left(\bar{\Omega}_{\nu}, \partial \Omega_{\nu}\right)$ yields

$$
\mathbb{R}=H^{0}\left(\bar{\Omega}_{\nu}, \mathbb{R}\right) \longrightarrow H^{0}\left(\partial \Omega_{\nu}, \mathbb{R}\right) \longrightarrow H^{1}\left(\bar{\Omega}_{\nu}, \partial \Omega_{\nu} ; \mathbb{R}\right) \simeq H_{c}^{1}\left(\Omega_{\nu}, \mathbb{R}\right)=0
$$

thus $H^{0}\left(\partial \Omega_{\nu}, \mathbb{R}\right)=\mathbb{R}$. Finally, the Mayer-Vietoris sequence applied to small neighborhoods of the two copies of $\bar{\Omega}_{\nu}$ in $Y_{\nu}$ gives an exact sequence

$$
H^{0}\left(\bar{\Omega}_{\nu}, \mathbb{R}\right)^{\oplus 2} \longrightarrow H^{0}\left(\partial \Omega_{\nu}, \mathbb{R}\right) \longrightarrow H^{1}\left(Y_{\nu}, \mathbb{R}\right) \longrightarrow H^{1}\left(\bar{\Omega}_{\nu}, \mathbb{R}\right)^{\oplus 2}=0
$$

where the first map is onto. Hence $H^{1}\left(Y_{\nu}, \mathbb{R}\right)=0$.
Proof End of the proof of the uniformization theorem.. Extend the almost complex structure of $\bar{\Omega}_{\nu}$ in an arbitrary way to $Y_{\nu}$, e.g. by an extension of a Riemannian metric. Then $Y_{\nu}$ becomes a compact Riemann surface of genus 0 , thus $Y_{\nu} \simeq \mathbb{P}^{1}$ and we obtain in particular a holomorphic embedding $\Phi_{\nu}: \Omega_{\nu} \rightarrow \mathbb{C}$. Fix a point $a \in \Omega_{0}$ and a non zero linear form $\xi^{\star} \in T_{a} X$. We can take the composition of $\Phi_{\nu}$ with an affine linear map $\mathbb{C} \rightarrow \mathbb{C}$ so that $\Phi_{\nu}(a)=0$ and $d \Phi_{\nu}(a)=\xi^{\star}$. By the well-known properties of injective holomorphic maps, $\left(\Phi_{\nu}\right)$ is then uniformly bounded on every small disk centered at $a$, thus also on every compact subset of $X$ by a connectedness argument. Hence ( $\Phi_{\nu}$ ) has a subsequence converging towards an injective holomorphic map $\Phi: X \rightarrow \mathbb{C}$.

## Chapter IX

## Finiteness Theorems for $q$-Convex Spaces and Stein Spaces

## § 1. Topological Preliminaries

## $\S$ 1.A. Krull Topology of $\mathfrak{O}_{\boldsymbol{n}}$-Modules

We shall use in an essential way different kind of topological results. The first of these concern the topology of modules over a local ring and depend on the Artin-Rees and Krull lemmas. Let $R$ be a noetherian local ring; "local" means that $R$ has a unique maximal ideal $\mathfrak{m}$, or equivalently, that $R$ has an ideal $\mathfrak{m}$ such that every element $\alpha \in R \backslash \mathfrak{m}$ is invertible.
(1.1) Nakayama lemma. Let $E$ be a finitely generated $R$-module such that $\mathfrak{m} E=E$. Then $E=\{0\}$.

Proof. By induction on the number of generators of $E$ : if $E$ is generated by $x_{1}, \ldots, x_{p}$, the hypothesis $E=\mathfrak{m} E$ shows that $x_{p}=\alpha_{1} x_{1}+\cdots+\alpha_{p} x_{p}$ with $\alpha_{j} \in \mathfrak{m}$; as $1-\alpha_{p} \in R \backslash \mathfrak{m}$ is invertible, we see that $x_{p}$ can be expressed in terms of $x_{1}, \ldots, x_{p-1}$ if $p>1$ and that $x_{1}=0$ if $p=1$.
(1.2) Artin-Rees lemma. Let $F$ be a finitely generated $R$-module and let $E$ be a submodule. There exists an integer s such that

$$
E \cap \mathfrak{m}^{k} F=\mathfrak{m}^{k-s}\left(E \cap \mathfrak{m}^{s} F\right) \quad \text { for } \quad k \geqslant s
$$

Proof. Let $R_{t}$ be the polynomial ring $R[\mathfrak{m} t]=R+\mathfrak{m} t+\cdots+\mathfrak{m}^{k} t^{k}+\cdots$ where $t$ is an indeterminate. If $g_{1}, \ldots, g_{p}$ is a set of generators of the ideal $\mathfrak{m}$ over $R$, we see that the ring $R_{t}$ is generated by $g_{1} t, \ldots, g_{p} t$ over $R$, hence $R_{t}$ is also noetherian. Now, we consider the $R_{t}$-modules

$$
E_{t}=\bigoplus E t^{k}, \quad F_{t}=\bigoplus\left(\mathfrak{m}^{k} F\right) t^{k}
$$

Then $F_{t}$ is generated over $R_{t}$ by the generators of $F$ over $R$, hence the submodule $E_{t} \cap F_{t}$ is finitely generated. Let $s$ be the highest exponent of $t$ in a set of generators $P_{1}(t), \ldots, P_{N}(t)$ of $E_{t} \cap F_{t}$. If we identify the components of $t^{k}$ in the extreme terms of the equality

$$
\bigoplus\left(E \cap \mathfrak{m}^{k} F\right) t^{k}=E_{t} \cap F_{t}=\sum_{j}\left(\bigoplus_{k} \mathfrak{m}^{k} t^{k}\right) P_{j}(t)
$$

we get

$$
E \cap \mathfrak{m}^{k} F \subset \sum_{l \leqslant s} \mathfrak{m}^{k-l}\left(E \cap \mathfrak{m}^{l} F\right) \subset \mathfrak{m}^{k-s}\left(E \cap \mathfrak{m}^{s} F\right)
$$

The opposite inclusion is clear.
(1.3) Krull lemma. Let $F$ be a finitely generated $R$-module and let $E$ be a submodule. Then
a) $\bigcap_{k \geqslant 0} \mathfrak{m}^{k} F=\{0\}$.
b) $\bigcap_{k \geqslant 0}\left(E+\mathfrak{m}^{k} F\right)=E$.

Proof. a) Put $G=\bigcap_{k \geqslant 0} \mathfrak{m}^{k} F \subset F$. By the Artin-Rees lemma, there exists $s \in \mathbb{N}$ such that $G \cap \mathfrak{m}^{k} F=\mathfrak{m}^{k-s}\left(G \cap \mathfrak{m}^{s} F\right)$. Taking $k=s+1$, we find $G \subset \mathfrak{m} G$, hence $\mathfrak{m} G=G$ and $G=\{0\}$ by the Nakayama lemma.
b) By applying a) to the quotient module $F / E$ we get $\cap \mathfrak{m}^{k}(F / E)=\{0\}$. Property b) follows.

Now assume that $R=\mathscr{O}_{n}=\mathbb{C}\left\{z_{1}, \ldots, z_{n}\right\}$ and $\mathfrak{m}=\left(z_{1}, \ldots, z_{n}\right)$. Then $\mathscr{O}_{n} / \mathfrak{m}^{k}$ is a finite dimensional vector space generated by the monomials $z^{\alpha},|\alpha|<k$. It follows that $E / \mathfrak{m}^{k} E$ is a finite dimensional vector space for any finitely generated $\mathscr{O}_{n}$-module $E$. As $\cap \mathfrak{m}^{k} E=\{0\}$ by 1.3 a), there is an injection

$$
\begin{equation*}
E \hookrightarrow \prod_{k \in \mathbb{N}} E / \mathfrak{m}^{k} E \tag{1.4}
\end{equation*}
$$

We endow $E$ with the Hausdorff topology induced by the product, i.e. with the weakest topology that makes all projections $E \longrightarrow E / \mathfrak{m}^{k} E$ continuous for the complex vector space topology on $E / \mathfrak{m}^{k} E$. This topology is called the Krull topology (or rather, the analytic Krull topology; the "algebraic" Krull topology would be obtained by taking the discrete topology on $\left.E / \mathfrak{m}^{k} E\right)$. For $E=\mathscr{O}_{n}$, this is the topology of simple convergence on coefficients, defined by the collection of semi-norms $\sum c_{\alpha} z^{\alpha} \longmapsto\left|c_{\alpha}\right|$. Observe that this topology is not complete: the completion of $\mathscr{Q}_{n}$ can be identified with the ring of formal power series $\mathbb{C}\left[\left[z_{1}, \ldots, z_{n}\right]\right]$. In general, the completion is the inverse limit $\widehat{E}=\lim E / \mathfrak{m}^{k} E$. Every $\widehat{O}_{n}$-homomorphism $E \longrightarrow F$ is continuous, because the induced

(1.5) Theorem. Let $E \subset F$ be finitely generated $\mathscr{O}_{n}$-modules. Then:
a) The map $F \longrightarrow G=F / E$ is open, i.e. the Krull topology of $G$ is the quotient of the Krull topology of $F$;
b) $E$ is closed in $F$ and the topology induced by $F$ on $E$ coincides with the Krull topology of $E$.

Proof. a) is an immediate consequence of the fact that the surjective finite dimensional linear maps $F / \mathfrak{m}^{k} F \longrightarrow G / \mathfrak{m}^{k} G$ are open.
b) Let $\bar{E}$ be the closure of $E$ in $F$. The image of $\bar{E}$ in $F / \mathfrak{m}^{k} F$ is mapped into the closure of the image of $E$. As every subspace of a finite dimensional space is closed, the images of $E$ and $\bar{E}$ must coincide, i.e. $\bar{E}+\mathfrak{m}^{k} F=E+\mathfrak{m}^{k} F$. Therefore

$$
E \subset \bar{E} \subset \bigcap\left(E+\mathfrak{m}^{k} F\right)=E
$$

thanks to 1.3 b ). The topology induced by $F$ on $E$ is the weakest that makes all projections $E \longrightarrow E / E \cap \mathfrak{m}^{k} F$ continuous (via the injections $E / E \cap \mathfrak{m}^{k} F \longleftrightarrow F / \mathfrak{m}^{k} F$ ). However, the Artin-Rees lemma gives

$$
\mathfrak{m}^{k} E \subset E \cap \mathfrak{m}^{k} F=\mathfrak{m}^{k-s}\left(E \cap \mathfrak{m}^{s} F\right) \subset \mathfrak{m}^{k-s} E \quad \text { for } \quad k \geqslant s
$$

so the topology induced by $F$ coincides with that induced by $\Pi E / \mathfrak{m}^{k} E$.

## $\S$ 1.B. Compact Perturbations of Linear Operators

We now recall some basic results in the perturbation theory of linear operators. These results will be needed in order to obtain a finiteness criterion for cohomology groups.
(1.6) Definition. Let $E, F$ be Hausdorff locally convex topological vector spaces and $g: E \longrightarrow F$ a continuous linear operator.
 $\overline{g(U)}$ is compact in $F$.
b) $g$ is said to be a monomorphism if $g$ is a topological isomorphism of $E$ onto a closed subspace of $F$, and a quasi-monomorphism if $\operatorname{ker} g$ is finite dimensional and $\tilde{g}: E / \operatorname{ker} g \longrightarrow F$ a monomorphism.
c) $g$ is said to be an epimorphism if $g$ is surjective and open, and a quasi-epimorphism if $g$ is an epimorphism of $E$ onto a closed finite codimensional subspace $F^{\prime} \subset F$.
d) $g$ is said to be a quasi-isomorphism if $g$ is simultaneously a quasi-monomorphism and a quasi-epimorphism.
(1.7) Lemma. Assume that $E, F$ are Fréchet spaces. Then
a) $g$ is a (quasi-) monomorphism if and only if $g(E)$ is closed in $F$ and $g$ is injective (resp. and ker $g$ is finite dimensional).
b) $g$ is a (quasi-) epimorphism if and only if $g$ is surjective (resp. $g(E)$ is finite codimensional).

Proof. a) If $g(E)$ is closed, the map $\widetilde{g}: E / \operatorname{ker} g \longrightarrow g(E)$ is a continuous bijective linear map between Fréchet spaces, so $\tilde{g}$ is a topological isomorphism by Banach's theorem.
b) If $g$ is surjective, Banach's theorem implies that $g$ is open, thus $g$ is an epimorphism. If $g(E)$ is finite codimensional, let $S$ be a supplementary subspace of $g(E)$ in $F, \operatorname{dim} S<$ $+\infty$. Then the map

$$
G: \quad(E / \operatorname{ker} g) \oplus S \longrightarrow F, \quad \widetilde{x} \oplus y \longmapsto \widetilde{g}(\widetilde{x})+y
$$

is a continuous bijective linear map between Fréchet spaces, so it is a topological isomorphism. In particular $g(E)=G((E / \operatorname{ker} g) \oplus\{0\})$ is closed as an image of a closed subspace. Hence $g(E)$ is also a Fréchet space and $g: E \longrightarrow g(E)$ is an epimorphism.
(1.8) Theorem (L. Schwartz). Let $h: E \longrightarrow F$ be a compact linear operator.
a) If $g: E \longrightarrow F$ is a quasi-monomorphism, then $g+h$ is a quasi-monomorphism.
b) If $E, F$ are Fréchet spaces and if $g: E \longrightarrow F$ is a quasi-epimorphism, then $g+h$ is a quasi-epimorphism.

Proof. Set $f=g+h$ and let $U$ be an open convex symmetric neighborhood of 0 in $E$ such that $K=\overline{h(U)}$ is compact.
a) It is sufficient to show that there is a closed finite codimensional subspace $E^{\prime} \subset E$ such that $f_{\left\lceil E^{\prime}\right.}$ is a monomorphism. If we first replace $E$ by a closed supplementary subspace of $\operatorname{ker} g$, we see that we may assume $g$ injective. Then $g$ is a monomorphism, so we may assume in fact that $E$ is a subspace of $F$ and that $g$ is the inclusion. Let $V$ be an open convex symmetric neighborhood of 0 in $F$ such that $V \cap E \subset U$. There exists a closed finite codimensional subspace $F^{\prime} \subset F$ such that $K \cap F^{\prime} \subset 2^{-1} V$ because $\bigcap_{F^{\prime}} K \cap F^{\prime}=\{0\}$. If we replace $E$ by $E^{\prime}=h^{-1}\left(F^{\prime}\right)$ and $U$ by $U^{\prime}=U \cap E^{\prime}$, we get

$$
K^{\prime}:=\overline{h\left(U^{\prime}\right)} \subset K \cap F^{\prime} \subset 2^{-1} V
$$

Hence, we may assume without loss of generality that $K \subset 2^{-1} V$. Then we show that $f=g+h$ is actually a monomorphism. If $\Omega$ is an arbitrary open neighborhood of 0 in $E$, we have to check that there exists a neighborhood $W$ of 0 in $F$ such that $f(x) \in W \Longrightarrow x \in \Omega$. There is an integer $N$ such that $2^{-N} K \cap E \subset \Omega$. We choose $W$ convex and so small that

$$
\left(W+2^{-N} K\right) \cap E \subset \Omega \quad \text { and } \quad 2^{N} W+K \subset 2^{-1} V
$$

Let $x \in E$ be such that $f(x) \in W$. Then $x \in 2^{n} U$ for some $n \in \mathbb{Z}$ (just take $n \geqslant n_{0}$ large enough). For any such $n \geqslant-N$ we infer

$$
x=g(x)=f(x)-h(x) \in W+2^{n} K=2^{n}\left(2^{-n} W+K\right) \subset 2^{n-1} V
$$

Thus $x \in 2^{n-1}(V \cap E) \subset 2^{n-1} U$. By induction we finally get $x \in 2^{-N} U$, so

$$
x \in\left(W+2^{-N} K\right) \cap E \subset \Omega .
$$

b) By Lemma 1.7 b ), we only have to show that there is a finite dimensional subspace $S \subset F$ such that the induced map

$$
\tilde{f}: E \longrightarrow F \longrightarrow F / S
$$

is surjective. If we take $S$ equal to a supplementary subspace of $g(E)$ and replace $g, h$ by the induced maps $\widetilde{g}, \widetilde{h}: E \longrightarrow F / S$, we may assume that $g$ itself is surjective. Then $g$ is open, so $V=g(U)$ is a convex open neighborhood of 0 in $F$. As $K$ is compact, there exists a finite set of elements $b_{1}, \ldots, b_{N} \in K$ such that $K \subset \bigcup\left(b_{j}+2^{-1} V\right)$. If we now
take $S=\operatorname{Vect}\left(b_{1}, \ldots, b_{N}\right)$ and replace again $g, h$ by $\widetilde{g}, \widetilde{h}$, we obtain $\widetilde{K} \subset 2^{-1} \widetilde{V}$ where $\widetilde{K}$ is the closure of $\widetilde{h}(U)$ and $V=\widetilde{g}(U)$, so we may assume in addition that $K \subset 2^{-1} V$. Then we show that $f=g+h$ is actually surjective. Let $y_{0} \in V$. There exists $x_{0} \in U$ such that $g\left(x_{0}\right)=y_{0}$, thus

$$
y_{1}=y_{0}-f\left(x_{0}\right)=-h\left(x_{0}\right) \in K \subset 2^{-1} V .
$$

By induction, we construct $x_{n} \in 2^{-n} U$ such that $g\left(x_{n}\right)=y_{n}$ and

$$
y_{n+1}=y_{n}-f\left(x_{n}\right)=-h\left(x_{n}\right) \in 2^{-n} K \subset 2^{-n-1} V .
$$

Hence $y_{n+1}=y_{0}-f\left(x_{0}+\cdots+x_{n}\right)$ tends to 0 in $F$ (as $y_{n+1} \in 2^{-n} K$ ), but we still have to make sure that the series $\sum x_{n}$ converges in $E$. Let $U_{p}$ be a fundamental system of convex neighborhoods of 0 in $E$ such that $U_{p+1} \subset 2^{-1} U_{p}$. For each $p, K$ is contained in the union of the open sets $g\left(2^{n} U_{p} \cap 2^{-1} U\right)$ when $n \in \mathbb{N}$, equal to $g\left(2^{-1} U\right)=2^{-1} V$. There exists an increasing sequence $N(p)$ such that $K \subset g\left(2^{N(p)} U_{p} \cap 2^{-1} U\right)$, thus

$$
2^{1-n} K \subset g\left(2^{N(p)+1-n} U_{p} \cap 2^{-n} U\right)
$$

As $y_{n} \in 2^{1-n} K$, we see that we can choose $x_{n} \in 2^{N(p)+1-n} U_{p} \cap 2^{-n} U$ for $N(p)<n \leqslant$ $N(p+1)$; then

$$
x_{N(p)+1}+\cdots+x_{N(p+1)} \in\left(1+2^{-1}+\cdots\right) U_{p} \subset 2 U_{p}
$$

As $E$ is complete, the series $x=\sum x_{n}$ converges towards an element $x \in E$ such that $f(x)=y_{0}$, and $f$ is surjective.

The following important finiteness theorem due to H. Cartan and J.-P. Serre can be easily deduced from this.
(1.9) Theorem. Let $\left(E^{\bullet}, d\right)$ and $\left(F^{\bullet}, \delta\right)$ be complexes of Fréchet spaces with continuous differentials, and $\rho^{\bullet}: E^{\bullet} \longrightarrow F^{\bullet}$ a continuous complex morphism. If $\rho^{q}$ is compact and $H^{q}\left(\rho^{\bullet}\right): H^{q}\left(E^{\bullet}\right) \longrightarrow H^{q}\left(F^{\bullet}\right)$ surjective, then $H^{q}\left(F^{\bullet}\right)$ is a Hausdorff finite dimensional space.

Proof. Consider the operators

$$
\begin{aligned}
& g, h: \quad Z^{q}\left(E^{\bullet}\right) \oplus F^{q-1} \longrightarrow Z^{q}\left(F^{\bullet}\right), \\
& g(x \oplus y)=\rho^{q}(x)+\delta^{q-1}(y), \quad h(x \oplus y)=-\rho^{q}(x) .
\end{aligned}
$$

As $Z^{q}\left(E^{\bullet}\right) \subset E^{q}, Z^{q}\left(F^{\bullet}\right) \subset F^{q}$ are closed, all our spaces are Fréchet spaces. Moreover the hypotheses imply that $h$ is compact and $g$ is surjective since $H^{q}\left(\rho^{\bullet}\right)$ is surjective. Hence $g$ is an epimorphism and $f=g+h=0 \oplus \delta^{q-1}$ is a quasi-epimorphism by 1.8 b ). Therefore $B^{q}\left(F^{\bullet}\right)$ is closed and finite codimensional in $Z^{q}\left(F^{\bullet}\right)$, thus $H^{q}\left(F^{\bullet}\right)$ is Hausdorff and finite dimensional.
(1.10) Remark. If $\rho^{\bullet}: E^{\bullet} \longrightarrow F^{\bullet}$ is a continuous morphism of Fréchet complexes and if $H^{q}\left(\rho^{\bullet}\right)$ is surjective, then $H^{q}\left(\rho^{\bullet}\right)$ is in fact open, because the above map $g$ is open. If $H^{q}\left(\rho^{\bullet}\right)$ is bijective, it follows that $H^{q}\left(\rho^{\bullet}\right)$ is necessarily a topological isomorphism (however $H^{q}\left(E^{\bullet}\right)$ and $H^{q}\left(F^{\bullet}\right)$ need not be Hausdorff).

## § 1.C. Abstract Mittag-Leffler Theorem

We will also need the following abstract Mittag-Leffler theorem, which is a very efficient tool in order to deal with cohomology groups of inverse limits.
(1.11) Proposition. Let $\left(E_{\nu}^{\bullet}, \delta\right)_{\nu \in \mathbb{N}}$ be a sequence of Fréchet complexes together with morphisms $E_{\nu+1}^{\bullet} \longrightarrow E_{\nu}^{\bullet}$. We assume that the image of $E_{\nu+1}^{\bullet}$ in $E_{\nu}^{\bullet}$ is dense and we let $E^{\bullet}=\lim _{\rightleftarrows} E_{\nu}^{\bullet}$ be the inverse limit complex.
a) If all maps $H^{q}\left(E_{\nu+1}^{\bullet}\right) \longrightarrow H^{q}\left(E_{\nu}^{\bullet}\right), \nu \in \mathbb{N}$, are surjective, then the limit $H^{q}\left(E^{\bullet}\right) \longrightarrow H^{q}\left(E_{0}^{\bullet}\right)$ is surjective.
b) If all maps $H^{q}\left(E_{\nu+1}^{\bullet}\right) \longrightarrow H^{q}\left(E_{\nu}^{\bullet}\right), \quad \nu \in \mathbb{N}$, have a dense range, then $H^{q}\left(E^{\bullet}\right) \longrightarrow H^{q}\left(E_{0}^{\bullet}\right)$ has a dense range.
c) If all maps $H^{q-1}\left(E_{\nu+1}^{\bullet}\right) \longrightarrow H^{q-1}\left(E_{\nu}^{\bullet}\right)$ have a dense range and all maps $H^{q}\left(E_{\nu+1}^{\bullet}\right) \longrightarrow H^{q}\left(E_{\nu}^{\bullet}\right)$ are injective, $\nu \in \mathbb{N}$, then $H^{q}\left(E^{\bullet}\right) \longrightarrow H^{q}\left(E_{0}^{\bullet}\right)$ is injective.
d) Let $\varphi^{\bullet}: F^{\bullet} \longrightarrow E^{\bullet}$ be a morphism of Fréchet complexes that has a dense range. If every map $H^{q}\left(F^{\bullet}\right) \longrightarrow H^{q}\left(E_{\nu}^{\bullet}\right)$ has a dense range, then $H^{q}\left(F^{\bullet}\right) \longrightarrow H^{q}\left(E^{\bullet}\right)$ has a dense range.

Proof. If $x$ is an element of $E^{\bullet}$ or of $E_{\mu}^{\bullet}, \mu \geqslant \nu$, we denote by $x^{\nu}$ its canonical image in $E_{\nu}^{\bullet}$. Let $d_{\nu}$ be a translation invariant distance that defines the topology of $E_{\nu}^{\bullet}$. After replacement of $d_{\nu}(x, y)$ by

$$
d_{\nu}^{\prime}(x, y)=\max _{\mu \leqslant \nu}\left\{d_{\mu}\left(x^{\mu}, y^{\mu}\right)\right\}, \quad x, y \in E_{\nu}^{\bullet}
$$

we may assume that all maps $E_{\nu+1}^{\bullet} \longrightarrow E_{\nu}^{\bullet}$ are Lipschitz continuous with coefficient 1.
a) Let $x_{0} \in Z^{q}\left(E_{0}^{\bullet}\right)$ represent a given cohomology class $\bar{x}_{0} \in H^{q}\left(E_{0}^{\bullet}\right)$. We construct by induction a convergent sequence $x_{\nu} \in Z^{q}\left(E_{\nu}^{\bullet}\right)$ such that $\bar{x}_{\nu}$ is mapped onto $\bar{x}_{0}$. If $x_{\nu}$ is already chosen, we can find by assumption $x_{\nu+1} \in Z^{q}\left(E_{\nu+1}^{\bullet}\right)$ such that $\bar{x}_{\nu+1}^{\nu}=\bar{x}_{\nu}$, i.e. $x_{\nu+1}^{\nu}=x_{\nu}+\delta y_{\nu}$ for some $y_{\nu} \in E_{\nu}^{q-1}$. If we replace $x_{\nu+1}$ by $x_{\nu+1}-\delta y_{\nu+1}$ where $y_{\nu+1} \in E_{\nu+1}^{q-1}$ yields an approximation $y_{\nu+1}^{\nu}$ of $y_{\nu}$, we may assume that

$$
\max \left\{d_{\nu}\left(y_{\nu}, 0\right), d_{\nu}\left(\delta y_{\nu}, 0\right)\right\} \leqslant 2^{-\nu}
$$

Then $\left(x_{\nu}\right)$ converges to a limit $\xi \in Z^{q}\left(E^{\bullet}\right)$ and we have $\xi^{0}=x_{0}+\delta \sum y_{\nu}^{0}$.
b) The density assumption for cohomology groups implies that the map

$$
Z^{q}\left(E_{\nu+1}^{\bullet}\right) \times E_{\nu}^{q-1} \longrightarrow Z^{q}\left(E_{\nu}^{\bullet}\right), \quad\left(x_{\nu+1}, y_{\nu}\right) \longmapsto x_{\nu+1}^{\nu}+\delta y_{\nu}
$$

has a dense range. If we approximate $y_{\nu}$ by elements coming from $E_{\nu+1}^{q-1}$, we see that the map $Z^{q}\left(E_{\nu+1}^{\bullet}\right) \longrightarrow Z^{q}\left(E_{\nu}^{\bullet}\right)$ has also a dense range. If $x_{0} \in Z^{q}\left(E_{0}^{\bullet}\right)$, we can find inductively a sequence $x_{\nu} \in Z^{q}\left(E_{\nu}^{\bullet}\right)$ such that $d_{\nu}\left(x_{\nu+1}^{\nu}, x_{\nu}\right) \leqslant \varepsilon 2^{-\nu-1}$ for all $\nu$, thus $\left(x_{\nu}\right)$ converges to an element $\xi \in Z^{q}\left(E^{\bullet}\right)$ such that $d_{0}\left(\xi^{0}, x_{0}\right) \leqslant \varepsilon$ and $Z^{q}\left(E^{\bullet}\right) \longrightarrow Z^{q}\left(E_{0}^{\bullet}\right)$ has a dense range.
c) Let $x \in Z^{q}\left(E^{\bullet}\right)$ be such that $\bar{x}^{0} \in H^{q}\left(E_{0}^{\bullet}\right)$ is zero. By assumption, the image of $\bar{x}$ in $H^{q}\left(E_{\nu}^{\bullet}\right)$ must be also zero, so we can write $x^{\nu}=d y_{\nu}, y_{\nu} \in E_{\nu}^{q-1}$. We have
$z_{\nu}=y_{\nu+1}^{\nu}-y_{\nu} \in Z^{q-1}\left(E_{\nu}^{\bullet}\right)$. Let $z_{\nu+1} \in Z^{q-1}\left(E_{\nu+1}^{\bullet}\right)$ be such that $z_{\nu+1}^{\nu}$ approximates $z_{\nu}$. If we replace $y_{\nu+1}$ by $y_{\nu+1}-z_{\nu+1}$, we still have $x^{\nu+1}=d y_{\nu+1}$ and we may assume in addition that $d_{\nu}\left(y_{\nu+1}^{\nu}, y_{\nu}\right) \leqslant 2^{-\nu}$. Then $\left(y_{\nu}\right)$ converges towards an element $y \in E^{q-1}$ such that $x=d y$, thus $\bar{x}=0$ and $H^{q}\left(E^{\bullet}\right) \longrightarrow H^{q}\left(E_{0}^{\bullet}\right)$ is injective.
d) For every class $\bar{y} \in H^{q}\left(E^{\bullet}\right)$, the hypothesis implies the existence of a sequence $x_{\nu} \in$ $Z^{q}\left(F^{\bullet}\right)$ such that $\varphi^{q}\left(\bar{x}_{\nu}\right)^{\nu}$ converges to $\bar{y}^{\nu}$, that is, $d_{\nu}\left(y^{\nu}, \varphi^{q}\left(x_{\nu}\right)^{\nu}+\delta z_{\nu}\right)$ tends to 0 for some sequence $z_{\nu} \in E_{\nu}^{q-1}$. Approximate $z_{\nu}$ by $\varphi^{q-1}\left(w_{\nu}\right)^{\nu}$ for some $w_{\nu} \in F^{q-1}$ and replace $x_{\nu}$ by $x_{\nu}^{\prime}=x_{\nu}+\delta w_{\nu}$. Then $\varphi^{q}\left(x_{\nu}^{\prime}\right)$ converges to $y$ in $Z^{q}\left(E^{\bullet}\right)$.

## §2. q-Convex Spaces

## $\S$ 2.A. $\boldsymbol{q}$-Convex Functions

The concept of $q$-convexity, first introduced in [Rothstein 1955], and further developed by [Andreotti-Grauert 1962], generalizes the concepts of pseudoconvexity already considered in chapters 1 and 8 . Let $M$ be a complex manifold, $\operatorname{dim}_{\mathbb{C}} M=n$. A function $v \in C^{2}(M, \mathbb{R})$ is said to be strongly (resp. weakly) $q$-convex at a point $x \in M$ if $i d^{\prime} d^{\prime \prime} v(x)$ has at least $(n-q+1)$ strictly positive (resp. nonnegative) eigenvalues, or equivalently if there exists a $(n-q+1)$-dimensional subspace $F \subset T_{x} M$ on which the complex Hessian $H_{x} v$ is positive definite (resp. semi-positive). Weak 1-convexity is thus equivalent to plurisubharmonicity. Some authors use different conventions for the number of positive eigenvalues in $q$-convexity. The reason why we introduce the number $n-q+1$ instead of $q$ is mainly due to the following result:
(2.1) Proposition. If $v \in C^{2}(M, \mathbb{R})$ is strongly (weakly) $q$-convex and if $Y$ is a submanifold of $M$, then $v_{\uparrow Y}$ is strongly (weakly) $q$-convex.

Proof. Let $d=\operatorname{dim} Y$. For every $x \in Y$, there exists $F \subset T_{x} M$ with $\operatorname{dim} F=n-q+1$ such that $H v$ is (semi-) positive on $F$. Then $G=F \cap T_{x} Y$ has dimension $\geqslant(n-q+1)-(n-d)=$ $d-q+1$, and $H\left(v_{\upharpoonright Y}\right)$ is (semi-) positive on $G \subset T_{x} Y$. Hence $v_{\upharpoonright Y}$ is strongly (weakly) $q$-convex at $x$.
(2.2) Proposition. Let $v_{j} \in C^{2}(M, \mathbb{R})$ be a weakly (strongly) $q_{j}$-convex function, $1 \leqslant$ $j \leqslant s$, and $\chi \in C^{2}\left(\mathbb{R}^{s}, \mathbb{R}\right)$ a convex function that is increasing (strictly increasing) in all variables. Then $v=\chi\left(v_{1}, \ldots, v_{s}\right)$ is weakly (strongly) $q$-convex with $q-1=\sum\left(q_{j}-1\right)$. In particular $v_{1}+\cdots+v_{s}$ is weakly (strongly) $q$-convex.

Proof. A simple computation gives

$$
\begin{equation*}
H v=\sum_{j} \frac{\partial \chi}{\partial t_{j}}\left(v_{1}, \ldots, v_{s}\right) H v_{j}+\sum_{j, k} \frac{\partial^{2} \chi}{\partial t_{j} \partial t_{k}}\left(v_{1}, \ldots, v_{s}\right) d^{\prime} v_{j} \otimes \overline{d^{\prime} v_{k}} \tag{2.3}
\end{equation*}
$$

and the second sum defines a semi-positive hermitian form. In every tangent space $T_{x} M$ there exists a subspace $F_{j}$ of codimension $q_{j}-1$ on which $H v_{j}$ is semi-positive (positive definite). Then $F=\bigcap F_{j}$ has codimension $\leqslant q-1$ and $H v$ is semi-positive (positive definite) on $F$.

The above result cannot be improved, as shown by the trivial example

$$
v_{1}(z)=-2\left|z_{1}\right|^{2}+\left|z_{2}\right|^{2}+\left|z_{3}\right|^{2}, \quad v_{2}(z)=\left|z_{1}\right|^{2}-2\left|z_{2}\right|^{2}+\left|z_{3}\right|^{2} \quad \text { on } \mathbb{C}^{3}
$$

in which case $q_{1}=q_{2}=2$ but $v_{1}+v_{2}$ is only 3 -convex. However, formula (2.3) implies the following result.
(2.4) Proposition. Let $v_{j} \in C^{2}(M, \mathbb{R}), 1 \leqslant j \leqslant s$, be such that every convex linear combination $\sum \alpha_{j} v_{j}, \alpha_{j} \geqslant 0, \sum \alpha_{j}=1$, is weakly (strongly) $q$-convex. If $\chi \in C^{2}\left(\mathbb{R}^{s}, \mathbb{R}\right)$ is a convex function that is increasing (strictly increasing) in all variables, then $\chi\left(v_{1}, \ldots, v_{s}\right)$ is weakly (strongly) $q$-convex.

The invariance property of Prop. 2.1 immediately suggests the definition of $q$-convexity on complex spaces or analytic schemes:
(2.5) Definition. Let $\left(X, \mathscr{O}_{X}\right)$ be an analytic scheme. A function $v$ on $X$ is said to be strongly (resp. weakly) $q$-convex of class $\mathscr{C}^{k}$ on $X$ if $X$ can be covered by patches $G: U \xrightarrow{\simeq} A, A \subset \Omega \subset \mathbb{C}^{N}$ such that for each patch there exists a function $\widetilde{v}$ on $\Omega$ with $\widetilde{v}_{\lceil A} \circ G=v_{\lceil U}$, which is strongly (resp. weakly) $q$-convex of class $\mathscr{C}^{k}$.

The notion of $q$-convexity on a patch $U$ does not depend on the way $U$ is embedded in $\mathbb{C}^{N}$, as shown by the following lemma.
(2.6) Lemma. Let $G: U \longrightarrow A \subset \Omega \subset \mathbb{C}^{N}$ and $G^{\prime}: U^{\prime} \longrightarrow A^{\prime} \subset \Omega^{\prime} \subset \mathbb{C}^{N^{\prime}}$ be two patches of $X$. Let $\widetilde{v}$ be a strongly (weakly) $q$-convex function on $\Omega$ and $v=\widetilde{v}_{\uparrow A} \circ G$. For every $x \in U \cap U^{\prime}$ there exists a strongly (weakly) $q$-convex function $\widetilde{v}^{\prime}$ on a neighborhood $W^{\prime} \subset \Omega^{\prime}$ of $G^{\prime}(x)$ such that $\widetilde{v}_{\uparrow A^{\prime} \cap W^{\prime}}^{\prime} \circ G^{\prime}$ coincides with $v$ on $G^{\prime-1}\left(W^{\prime}\right)$.

Proof. The isomorphisms

$$
\begin{aligned}
& G^{\prime} \circ G^{-1}: A \supset G\left(U \cap U^{\prime}\right) \longrightarrow G^{\prime}\left(U \cap U^{\prime}\right) \subset A^{\prime} \\
& G \circ G^{\prime-1}: A^{\prime} \supset G^{\prime}\left(U \cap U^{\prime}\right) \longrightarrow G\left(U \cap U^{\prime}\right) \subset A
\end{aligned}
$$

are restrictions of holomorphic maps $H: W \longrightarrow \Omega^{\prime}, H^{\prime}: W^{\prime} \longrightarrow \Omega$ defined on neighborhoods $W \ni G(x), W^{\prime} \ni G^{\prime}(x)$; we can shrink $W^{\prime}$ so that $H^{\prime}\left(W^{\prime}\right) \subset W$. If we compose the automorphism $\left(z, z^{\prime}\right) \longmapsto\left(z, z^{\prime}-H(z)\right)$ of $W \times \mathbb{C}^{N^{\prime}}$ with the function $v(z)+\left|z^{\prime}\right|^{2}$ we see that the function $\varphi\left(z, z^{\prime}\right)=\widetilde{v}(z)+\left|z^{\prime}-H(z)\right|^{2}$ is strongly (weakly) $q$-convex on $W \times \Omega^{\prime}$. Now, $W^{\prime}$ can be embedded in $W \times \Omega^{\prime}$ via the map $z^{\prime} \longmapsto\left(H^{\prime}\left(z^{\prime}\right), z^{\prime}\right)$, so that the composite function

$$
\widetilde{v}^{\prime}\left(z^{\prime}\right)=\varphi\left(H^{\prime}\left(z^{\prime}\right), z^{\prime}\right)=\widetilde{v}\left(H^{\prime}\left(z^{\prime}\right)\right)+\left|z^{\prime}-H \circ H^{\prime}\left(z^{\prime}\right)\right|^{2}
$$

is strongly (weakly) $q$-convex on $W^{\prime}$ by Prop. 2.1. Since $H \circ G=G^{\prime}$ and $H^{\prime} \circ G^{\prime}=G$ on $G^{\prime-1}\left(W^{\prime}\right)$, we have $\widetilde{v}^{\prime} \circ G^{\prime}=\widetilde{v} \circ G=v$ on $G^{\prime-1}\left(W^{\prime}\right)$ and the lemma follows.

A consequence of this lemma is that Prop. 2.2 is still valid for an analytic scheme $X$ (all the extensions $\widetilde{v}_{j}$ near a given point $x \in X$ can be obtained with respect to the same local embedding).
(2.7) Definition. An analytic scheme $\left(X, \mathscr{O}_{X}\right)$ is said to be strongly (resp. weakly) $q$ convex if $X$ has a $\mathscr{C}^{\infty}$ exhaustion function $\psi$ which is strongly (resp. weakly) $q$-convex outside an exceptional compact set $K \subset X$. We say that $X$ is strongly $q$-complete if $\psi$
can be chosen so that $K=\emptyset$. By convention, a compact scheme $X$ is said to be strongly 0 -complete, with exceptional compact set $K=X$.

We consider the sublevel sets

$$
\begin{equation*}
X_{c}=\{x \in X ; \psi(x)<c\}, \quad c \in \mathbb{R} \tag{2.8}
\end{equation*}
$$

If $K \subset X_{c}$, we may select a convex increasing function $\chi$ such that $\chi=0$ on $\left.]-\infty, c\right]$ and $\chi^{\prime}>0$ on $] c,+\infty\left[\right.$. Then $\chi \circ \psi=0$ on $X_{c}$, so that $\chi \circ \psi$ is weakly $q$-convex everywhere in virtue of (2.3). In the weakly $q$-convex case, we may therefore always assume $K=\emptyset$. The following properties are almost immediate consequences of the definition:

## (2.9) Theorem.

a) A scheme $X$ is strongly (weakly) q-convex if and only if the reduced space $X_{\text {red }}$ is strongly (weakly) $q$-convex.
b) If $X$ is strongly (weakly) q-convex, every closed analytic subset $Y$ of $X_{\text {red }}$ is strongly (weakly) $q$-convex.
c) If $X$ is strongly (weakly) q-convex, every sublevel set $X_{c}$ containing the exceptional compact set $K$ is strongly (weakly) $q$-convex.
d) If $U_{j}$ is a weakly $q_{j}$-convex open subset of $X, 1 \leqslant j \leqslant s$, the intersection $U=$ $U_{1} \cap \ldots \cap U_{s}$ is weakly $q$-convex with $q-1=\sum\left(q_{j}-1\right) ; U$ is strongly $q$-convex (resp. $q$-complete) as soon as one of the sets $U_{j}$ is strongly $q_{j}$-convex (resp. $q_{j}$-complete).

Proof. a) is clear, since Def. 2.5 does not involve the structure sheaf $\mathbb{O}_{X}$. In cases b) and c), let $\psi$ be an exhaustion of the required type on $X$. Then $\psi_{Y Y}$ and $1 /(c-\psi)$ are exhaustions on $Y$ and $X_{c}$ respectively (this is so only if $Y$ is closed). Moreover, these functions are strongly (weakly) $q$-convex on $Y \backslash(K \cap Y)$ and $X_{c} \backslash K$, thanks to Prop. 2.1 and 2.2. For property d), note that a sum $\psi=\psi_{1}+\cdots+\psi_{s}$ of exhaustion functions on the sets $U_{j}$ is an exhaustion on $U$, choose the $\psi_{j}$ 's weakly $q_{j}$-convex everywhere, and apply Prop. 2.2.
(2.10) Corollary. Any finite intersection $U=U_{1} \cap \ldots \cap U_{s}$ of weakly 1-convex open subsets is weakly 1-convex. The set $U$ is strongly 1-convex (resp. 1-complete) as soon as one of the sets $U_{j}$ is strongly 1-convex (resp. 1-complete).

## § 2.B. Neighborhoods of $\boldsymbol{q}$-complete subspaces

We prove now a rather useful result asserting the existence of $q$-complete neighborhoods for $q$-complete subvarieties. The case $q=1$ goes back to [Siu 1976], who used a much more complicated method. The first step is an approximation-extension theorem for strongly $q$-convex functions.
(2.11) Proposition. Let $Y$ be an analytic set in a complex space $X$ and $\psi$ a strongly $q$-convex $\mathscr{C}^{\infty}$ function on $Y$. For every continuous function $\delta>0$ on $Y$, there exists a strongly $q$-convex $\mathscr{C}^{\infty}$ function $\varphi$ on a neighborhood $V$ of $Y$ such that $\psi \leqslant \varphi_{\mid Y}<\psi+\delta$.

Proof. Let $Z_{k}$ be a stratication of $Y$ as given by Prop. II.5.6, i.e. $Z_{k}$ is an increasing sequence of analytic subsets of $Y$ such that $Y=\bigcup Z_{k}$ and $Z_{k} \backslash Z_{k-1}$ is a smooth $k$ dimensional manifold (possibly empty for some $k$ 's). We shall prove by induction on $k$ the following statement:

There exists a $\mathscr{C}^{\infty}$ function $\varphi_{k}$ on $X$ which is strongly $q$-convex along $Y$ and on a closed neighborhood $\bar{V}_{k}$ of $Z_{k}$ in $X$, such that $\psi \leqslant \varphi_{k \upharpoonright Y}<\psi+\delta$.

We first observe that any smooth extension $\varphi_{-1}$ of $\psi$ to $X$ satisfies the requirements with $Z_{-1}=V_{-1}=\emptyset$. Assume that $V_{k-1}$ and $\varphi_{k-1}$ have been constructed. Then $Z_{k} \backslash V_{k-1} \subset Z_{k} \backslash Z_{k-1}$ is contained in $Z_{k, \text { reg }}$. The closed set $Z_{k} \backslash V_{k-1}$ has a locally finite covering $\left(A_{\lambda}\right)$ in $X$ by open coordinate patches $A_{\lambda} \subset \Omega_{\lambda} \subset \mathbb{C}^{N_{\lambda}}$ in which $Z_{k}$ is given by equations $z_{\lambda}^{\prime}=\left(z_{\lambda, k+1}, \ldots, z_{\lambda, N_{\lambda}}\right)=0$. Let $\theta_{\lambda}$ be $\mathscr{C}^{\infty}$ functions with compact support in $A_{\lambda}$ such that $0 \leqslant \theta_{\lambda} \leqslant 1$ and $\sum \theta_{\lambda}=1$ on $Z_{k} \backslash V_{k-1}$. We set

$$
\varphi_{k}(x)=\varphi_{k-1}(x)+\sum \theta_{\lambda}(x) \varepsilon_{\lambda}^{3} \log \left(1+\varepsilon_{\lambda}^{-4}\left|z_{\lambda}^{\prime}\right|^{2}\right) \quad \text { on } \quad X .
$$

For $\varepsilon_{\lambda}>0$ small enough, we will have $\psi \leqslant \varphi_{k-1 \upharpoonright Y} \leqslant \varphi_{k \upharpoonright Y}<\psi+\delta$. Now, we check that $\varphi_{k}$ is still strongly $q$-convex along $Y$ and near any $x_{0} \in \bar{V}_{k-1}$, and that $\varphi_{k}$ becomes strongly $q$-convex near any $x_{0} \in Z_{k} \backslash V_{k-1}$. We may assume that $x_{0} \in \operatorname{Supp} \theta_{\mu}$ for some $\mu$, otherwise $\varphi_{k}$ coincides with $\varphi_{k-1}$ in a neighborhood of $x_{0}$. Select $\mu$ and a small neighborhood $W \subset \subset \Omega_{\mu}$ of $x_{0}$ such that
a) if $x_{0} \in Z_{k} \backslash V_{k-1}$, then $\theta_{\mu}\left(x_{0}\right)>0$ and $A_{\mu} \cap W \subset \subset\left\{\theta_{\mu}>0\right\}$;
b) if $x_{0} \in A_{\lambda}$ for some $\lambda$ (there is only a finite set $I$ of such $\lambda$ 's), then $A_{\mu} \cap W \subset \subset A_{\lambda}$ and $z_{\lambda \mid A_{\mu} \cap W}$ has a holomorphic extension $\widetilde{z}_{\lambda}$ to $\bar{W}$;
c) if $x_{0} \in \bar{V}_{k-1}$, then $\varphi_{k-1 \upharpoonright A_{\mu} \cap W}$ has a strongly $q$-convex extension $\widetilde{\varphi}_{k-1}$ to $\bar{W}$;
d) if $x_{0} \in Y \backslash \bar{V}_{k-1}$, then $\varphi_{k-1 \upharpoonright Y \cap W}$ has a strongly $q$-convex extension $\widetilde{\varphi}_{k-1}$ to $\bar{W}$.

Otherwise take an arbitrary smooth extension $\widetilde{\varphi}_{k-1}$ of $\varphi_{k-1 \upharpoonright A_{\mu} \cap W}$ to $\bar{W}$ and let $\widetilde{\theta}_{\lambda}$ be an extension of $\theta_{\lambda\left\lceil A_{\mu} \cap W\right.}$ to $\bar{W}$. Then

$$
\widetilde{\varphi}_{k}=\widetilde{\varphi}_{k-1}+\sum \widetilde{\theta}_{\lambda} \varepsilon_{\lambda}^{3} \log \left(1+\varepsilon_{\lambda}^{-4}\left|\widetilde{\widetilde{z}}_{\lambda}^{\prime}\right|^{2}\right)
$$

is an extension of $\varphi_{k \upharpoonright A_{\mu} \cap W}$ to $\bar{W}$, resp. of $\varphi_{k \upharpoonright Y \cap W}$ to $\bar{W}$ in case d). As the function $\log \left(1+\varepsilon_{\lambda}^{-4}\left|\widetilde{z}_{\lambda}^{\prime}\right|^{2}\right)$ is plurisubharmonic and as its first derivative $\left\langle\widetilde{z}_{\lambda}^{\prime}, d \widetilde{z}_{\lambda}^{\prime}\right\rangle\left(\varepsilon_{\lambda}^{4}+\left|\widetilde{z}_{\lambda}^{\prime}\right|^{2}\right)^{-1}$ is bounded by $O\left(\varepsilon_{\lambda}^{-2}\right)$, we see that

$$
i d^{\prime} d^{\prime \prime} \widetilde{\varphi}_{k} \geqslant i d^{\prime} d^{\prime \prime} \widetilde{\varphi}_{k-1}-O\left(\sum \varepsilon_{\lambda}\right)
$$

Therefore, for $\varepsilon_{\lambda}$ small enough, $\widetilde{\varphi}_{k}$ remains $q$-convex on $\bar{W}$ in cases c) and d). Since all functions $\widetilde{z}_{\lambda}^{\prime}$ vanish along $Z_{k} \cap W$, we have

$$
i d^{\prime} d^{\prime \prime} \widetilde{\varphi}_{k} \geqslant i d^{\prime} d^{\prime \prime} \widetilde{\varphi}_{k-1}+\sum_{\lambda \in I} \theta_{\lambda} \varepsilon_{\lambda}^{-1} i d^{\prime} d^{\prime \prime}\left|\widetilde{z}_{\lambda}^{\prime}\right|^{2} \geqslant i d^{\prime} d^{\prime \prime} \widetilde{\varphi}_{k-1}+\theta_{\mu} \varepsilon_{\mu}^{-1} i d^{\prime} d^{\prime \prime}\left|z_{\mu}^{\prime}\right|^{2}
$$

at every point of $Z_{k} \cap W$. Moreover $i d^{\prime} d^{\prime \prime} \widetilde{\varphi}_{k-1}$ has at most ( $q-1$ )-negative eigenvalues on $T Z_{k}$ since $Z_{k} \subset Y$, whereas $i d^{\prime} d^{\prime \prime}\left|z_{\mu}^{\prime}\right|^{2}$ is positive definite in the normal directions
to $Z_{k}$ in $\Omega_{\mu}$. In case a), we thus find that $\widetilde{\varphi}_{k}$ is strongly $q$-convex on $\bar{W}$ for $\varepsilon_{\mu}$ small enough; we also observe that only finitely many conditions are required on each $\varepsilon_{\lambda}$ if we choose a locally finite covering of $\bigcup$ Supp $\theta_{\lambda}$ by neighborhoods $W$ as above. Therefore, for $\varepsilon_{\lambda}$ small enough, $\varphi_{k}$ is strongly $q$-convex on a neighborhood $\bar{V}_{k}^{\prime}$ of $Z_{k} \backslash V_{k-1}$. The function $\varphi_{k}$ and the set $V_{k}=V_{k}^{\prime} \cup V_{k-1}$ satisfy the requirements at order $k$. It is clear that we can choose the sequence $\varphi_{k}$ stationary on every compact subset of $X$; the limit $\varphi$ and the open set $V=\bigcup V_{k}$ fulfill the proposition.

The second step is the existence of almost plurisubharmonic functions having poles along a prescribed analytic set. By an almost plurisubharmonic function on a manifold, we mean a function that is locally equal to the sum of a plurisubharmonic function and of a smooth function, or equivalently, a function whose complex Hessian has bounded negative part. On a complex space, we require that our function can be locally extended as an almost plurisubharmonic function in the ambient space of an embedding.
(2.12) Lemma. Let $Y$ be an analytic subvariety in a complex space $X$. There is an almost plurisubharmonic function $v$ on $X$ such that $v=-\infty$ on $Y$ with logarithmic poles and $v \in \mathscr{C}^{\infty}(X \backslash Y)$.

Proof. Since $\mathscr{I}_{Y} \subset \mathscr{O}_{X}$ is a coherent subsheaf, there is a locally finite covering of $X$ by patches $A_{\lambda}$ isomorphic to analytic sets in balls $B\left(0, r_{\lambda}\right) \subset \mathbb{C}^{N_{\lambda}}$, such that $\mathscr{I}_{Y}$ admits a system of generators $g_{\lambda}=\left(g_{\lambda, j}\right)$ on a neighborhood of each set $\bar{A}_{\lambda}$. We set

$$
\begin{aligned}
& v_{\lambda}(z)=\log \left|g_{\lambda}(z)\right|^{2}-\frac{1}{r_{\lambda}^{2}-\left|z-z_{\lambda}\right|^{2}} \quad \text { on } A_{\lambda} \\
& v(z)=M_{(1, \ldots, 1)}\left(\ldots, v_{\lambda}(z), \ldots\right) \quad \text { for } \lambda \text { such that } A_{\lambda} \ni z
\end{aligned}
$$

where $M_{\eta}$ is the regularized max function defined in I-3.37. As the generators $\left(g_{\lambda, j}\right)$ can be expressed in terms of one another on a neighborhood of $\bar{A}_{\lambda} \cap \bar{A}_{\mu}$, we see that the quotient $\left|g_{\lambda}\right| /\left|g_{\mu}\right|$ remains bounded on this set. Therefore none of the values $v_{\lambda}(z)$ for $A_{\lambda} \ni z$ and $z$ near $\partial A_{\lambda}$ contributes to the value of $v$, since $1 /\left(r_{\lambda}^{2}-\left|z-z_{\lambda}\right|^{2}\right)$ tends to $+\infty$ on $\partial A_{\lambda}$. It follows that $v$ is smooth on $X \backslash Y$; as each $v_{\lambda}$ is almost plurisubharmonic on $A_{\lambda}$, we also see that $v$ is almost plurisubharmonic on $X$.
(2.13) Theorem. Let $X$ be a complex space and $Y$ a strongly $q$-complete analytic subset. Then $Y$ has a fundamental family of strongly $q$-complete neighborhoods $V$ in $X$.

Proof. By Prop. 2.11 applied to a strongly $q$-convex exhaustion of $Y$ and $\delta=1$, there exists a strongly $q$-convex function $\varphi$ on a neighborhood $W_{0}$ of $Y$ such that $\varphi_{\upharpoonright Y}$ is an exhaustion. Let $W_{1}$ be a neighborhood of $Y$ such that $\bar{W}_{1} \subset W_{0}$ and such that $\varphi_{\upharpoonright} \bar{W}_{1}$ is an exhaustion. We are going to show that every neighborhood $W \subset W_{1}$ of $Y$ contains a strongly $q$-complete neighborhood $V$. If $v$ is the function given by Lemma 2.12, we set

$$
\widetilde{v}=v+\chi \circ \varphi \quad \text { on } \bar{W}
$$

where $\chi: \mathbb{R} \rightarrow \mathbb{R}$ is a smooth convex increasing function. If $\chi$ grows fast enough, we get $\widetilde{v}>0$ on $\partial W$ and the $(q-1)$-codimensional subspace on which $i d^{\prime} d^{\prime \prime} \varphi$ is positive definite (in some ambient space) is also positive definite for $i d^{\prime} d^{\prime \prime} \widetilde{v}$ provided that $\chi^{\prime}$ be large enough to compensate the bounded negative part of $i d^{\prime} d^{\prime \prime} v$. Then $\widetilde{v}$ is strongly $q$-convex. Let $\theta$ be a smooth convex increasing function on $]-\infty, 0[$ such that $\theta(t)=0$ for $t<-3$ and $\theta(t)=-1 / t$ on $]-1,0[$. The open set $V=\{z \in W ; \widetilde{v}(z)<0\}$ is a neighborhood of $Y$ and $\psi=\varphi+\theta \circ \widetilde{v}$ is a strongly $q$-convex exhaustion of $V$.

## § 2.C. Runge Open Subsets

In order to extend the classical Runge theorem into an approximation result for sheaf cohomology groups, we need the concept of a $q$-Runge open subset.
(2.14) Definition. An open subset $U$ of a complex space $X$ is said to be $q$-Runge (resp. $q$-Runge complete) in $X$ if for every compact subset $L \subset U$ there exists a smooth exhaustion function $\psi$ on $X$ and a sublevel set $X_{b}$ of $\psi$ such that $L \subset X_{b} \subset \subset U$ and $\psi$ is strongly $q$-convex on $X \backslash \bar{X}_{b}$ (resp. on the whole space $X$ ).
(2.15) Example. If $X$ is strongly $q$-complete and if $\psi$ is a strongly $q$-convex exhaustion function of $X$, then every sublevel set $X_{c}$ of $\psi$ is $q$-Runge complete in $X$ : every compact set $L \subset X_{c}$ satisfies $L \subset X_{b} \subset \subset X_{c}$ for some $b<c$. More generally, if $X$ is strongly $q$-convex and if $\psi$ is strongly $q$-convex on $X \backslash K$, every sublevel set $X_{c}$ containing $K$ is $q$-Runge in $X$.

Later on, we shall need the following technical result.
(2.16) Proposition. Let $Y$ be an analytic subset of a complex space $X$. If $U$ is a $q$-Runge complete open subset of $Y$ and $L$ a compact subset, there exist a neighborhood $V$ of $Y$ in $X$ and a strongly $q$-convex exhaustion $\widetilde{\psi}$ on $V$ such that $U=Y \cap V$ and $L \subset Y \cap V_{b} \subset \subset U$ for some sublevel set $V_{b}$ of $\widetilde{\psi}$.

Proof. Let $\psi$ be a strongly $q$-convex exhaustion on $Y$ with $L \subset\{\psi<b\} \subset \subset U$ as in Def. 2.14. Then $L \subset\{\psi<b-\delta\}$ for some number $\delta>0$ and Lemma 2.11 gives a strongly $q$-convex function $\varphi$ on a neighborhood $W_{0}$ of $Y$ so that $\psi \leqslant \varphi_{\upharpoonright} Y<\psi+\delta$. The neighborhood $V$ and the function $\widetilde{\psi}=\varphi+\theta \circ \widetilde{v}$ constructed in the proof of Th. 2.13 are the desired ones: we have $\psi \leqslant \widetilde{\psi}_{\Gamma Y}=\varphi_{\mid Y}<\psi+\delta$, thus

$$
L \subset Y \cap V_{b-\delta} \subset\{\psi<b\} \subset \subset U
$$

## § 3. $\boldsymbol{q}$-Convexity Properties in Top Degrees

It is obvious by definition that a $n$-dimensional complex manifold $M$ is strongly $q$ complete for $q \geqslant n+1$ (an arbitrary smooth function is then strongly $q$-convex !). If $M$ is connected and non compact, [Greene and Wu 1975] have shown that $M$ is strongly $n$-complete, i.e. there is a smooth exhaustion function $\psi$ on $M$ such that $i d^{\prime} d^{\prime \prime} \psi$ has at least one positive eigenvalue everywhere. We need the following lemmas.
(3.1) Lemma. Let $\psi$ be a strongly $q$-convex function on $M$ and $\varepsilon>0$ a given number. There exists a hermitian metric $\omega$ on $M$ such that the eigenvalues $\gamma_{1} \leqslant \ldots \leqslant \gamma_{n}$ of the Hessian form $i d^{\prime} d^{\prime \prime} \psi$ with respect to $\omega$ satisfy $\gamma_{1} \geqslant-\varepsilon$ and $\gamma_{q}=\ldots=\gamma_{n}=1$.

Proof. Let $\omega_{0}$ be a fixed hermitian metric, $A_{0} \in \mathscr{C}^{\infty}(\operatorname{End} T M)$ the hermitian endomorphism associated to the hermitian form $i d^{\prime} d^{\prime \prime} \psi$ with respect to $\omega_{0}$, and $\gamma_{1}^{0} \leqslant \ldots \leqslant \gamma_{n}^{0}$ the eigenvalues of $A_{0}$ (or $i d^{\prime} d^{\prime \prime} \psi$ ). We can choose a function $\eta \in \mathscr{C}^{\infty}(M, \mathbb{R})$ such that $0<\eta(x) \leqslant \gamma_{q}^{0}(x)$ at each point $x \in M$. Select a positive function $\theta \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$ such that

$$
\theta(t) \geqslant|t| / \varepsilon \text { for } t \leqslant 0, \quad \theta(t) \geqslant t \text { for } t \geqslant 0, \quad \theta(t)=t \text { for } t \geqslant 1
$$

We let $\omega$ be the hermitian metric defined by the hermitian endomorphism

$$
A(x)=\eta(x) \theta\left[(\eta(x))^{-1} A_{0}(x)\right]
$$

where $\theta\left[\eta^{-1} A_{0}\right] \in \mathscr{C}^{\infty}(\operatorname{End} T M)$ is defined as in Lemma VII-6.2. By construction, the eigenvalues of $A(x)$ are $\alpha_{j}(x)=\eta(x) \theta\left(\gamma_{j}^{0}(x) / \eta(x)\right)>0$ and we have

$$
\begin{array}{ll}
\alpha_{j}(x) \geqslant\left|\gamma_{j}^{0}(x)\right| / \varepsilon & \text { for } \gamma_{j}^{0}(x) \leqslant 0 \\
\alpha_{j}(x) \geqslant \gamma_{j}^{0}(x) & \text { for } \gamma_{j}^{0}(x) \geqslant 0, \\
\alpha_{j}(x)=\gamma_{j}^{0}(x) & \text { for } j \geqslant q \quad\left(\text { then } \gamma_{j}^{0}(x) \geqslant \eta(x)\right)
\end{array}
$$

The eigenvalues of $i d^{\prime} d^{\prime \prime} \psi$ with respect to $\omega$ are $\gamma_{j}(x)=\gamma_{j}^{0}(x) / \alpha_{j}(x)$ and they have the required properties.

On a hermitian manifold $(M, \omega)$, we consider the Laplace operator $\Delta_{\omega}$ defined by

$$
\begin{equation*}
\Delta_{\omega} v=\operatorname{Trace}_{\omega}\left(i d^{\prime} d^{\prime \prime} v\right)=\sum_{1 \leqslant j, k \leqslant n} \omega^{j k}(z) \frac{\partial^{2} v}{\partial z_{j} \partial \bar{z}_{k}} \tag{3.2}
\end{equation*}
$$

where $\left(\omega^{j k}\right)$ is the conjugate of the inverse matrix of $\left(\omega_{j k}\right)$. Note that $\Delta_{\omega}$ may differ from the usual Laplace-Beltrami operator if $\omega$ is not Kähler. We say that $v$ is strongly $\omega$-subharmonic if $\Delta_{\omega} v>0$. This property implies clearly that $v$ is strongly $n$-convex; however, as

$$
\begin{aligned}
\Delta_{\omega} \chi\left(v_{1}, \ldots, v_{s}\right)=\sum_{j} & \frac{\partial \chi}{\partial t_{j}}\left(v_{1}, \ldots, v_{s}\right) \Delta_{\omega} v_{j} \\
& +\sum_{j, k} \frac{\partial^{2} \chi}{\partial t_{j} \partial t_{k}}\left(v_{1}, \ldots, v_{s}\right)\left\langle d^{\prime} v_{j}, d^{\prime} v_{k}\right\rangle_{\omega}
\end{aligned}
$$

subharmonicity has the advantage of being preserved by all convex increasing transformations. Conversely, if $\psi$ is strongly $n$-convex and $\omega$ chosen as in Lemma 3.1 with $\varepsilon$ small enough, we get $\Delta_{\omega} \psi \geqslant 1-(n-1) \varepsilon>0$, thus $\psi$ is strongly subharmonic for a suitable metric $\omega$.
(3.3) Lemma. Let $U, W \subset M$ be open sets such that for every connected component $U_{s}$ of $U$ there is a connected component $W_{t(s)}$ of $W$ such that $W_{t(s)} \cap U_{s} \neq \emptyset$ and $W_{t(s)} \backslash \bar{U}_{s} \neq \emptyset$. Then there exists a function $v \in \mathscr{C}^{\infty}(M, \mathbb{R}), v \geqslant 0$, with support contained in $\bar{U} \cup \bar{W}$, such that $v$ is strongly $\omega$-subharmonic and $>0$ on $U$.

Proof. We first prove that the result is true when $U, W$ are small cylinders with the same radius and axis. Let $a_{0} \in M$ be a given point and $z_{1}, \ldots, z_{n}$ holomorphic coordinates centered at $a_{0}$. We set $\operatorname{Re} z_{j}=x_{2 j-1}, \operatorname{Im} z_{j}=x_{2 j}, x^{\prime}=\left(x_{2}, \ldots, x_{2 n}\right)$ and $\omega=\sum \widetilde{\omega}_{j k}(x) d x_{j} \otimes d x_{k}$. Let $U$ be the cylinder $\left|x_{1}\right|<r,\left|x^{\prime}\right|<r$, and $W$ the cylinder $r-\varepsilon<x_{1}<r+\varepsilon,\left|x^{\prime}\right|<r$. There are constants $c, C>0$ such that

$$
\sum \widetilde{\omega}^{j k}(x) \xi_{j} \xi_{k} \geqslant c|\xi|^{2} \quad \text { and } \quad \sum\left|\widetilde{\omega}^{j k}(x)\right| \leqslant C \quad \text { on } \bar{U}
$$

Let $\chi \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$ be a nonnegative function equal to 0 on $\left.]-\infty,-r\right] \cup[r+\varepsilon,+\infty[$ and strictly convex on ] -r,r]. We take explicitly $\chi\left(x_{1}\right)=\left(x_{1}+r\right) \exp \left(-1 /\left(x_{1}+r\right)^{2}\right)$ on ] $-r, r]$ and

$$
v(x)=\chi\left(x_{1}\right) \exp \left(1 /\left(\left|x^{\prime}\right|^{2}-r^{2}\right)\right) \quad \text { on } \quad U \cup W, \quad v=0 \quad \text { on } \quad M \backslash(U \cup W)
$$

We have $v \in C^{\infty}(M, \mathbb{R}), v>0$ on $U$, and a simple computation gives

$$
\begin{aligned}
& \frac{\Delta_{\omega} v(x)}{v(x)}=\widetilde{\omega}^{11}(x)\left(4\left(x_{1}+r\right)^{-5}-2\left(x_{1}+r\right)^{-3}\right) \\
& \quad+\sum_{j>1} \widetilde{\omega}^{1 j}(x)\left(1+2\left(x_{1}+r\right)^{-2}\right)\left(-2 x_{j}\right)\left(r^{2}-\left|x^{\prime}\right|^{2}\right)^{-2} \\
& \quad+\sum_{j, k>1} \widetilde{\omega}^{j k}(x)\left(x_{j} x_{k}\left(4-8\left(r^{2}-\left|x^{\prime}\right|^{2}\right)\right)-2\left(r^{2}-\left|x^{\prime}\right|^{2}\right)^{2} \delta_{j k}\right)\left(r^{2}-\left|x^{\prime}\right|^{2}\right)^{-4}
\end{aligned}
$$

For $r$ small, we get

$$
\begin{aligned}
\frac{\Delta_{\omega} v(x)}{v(x)} \geqslant 2 c\left(x_{1}+r\right)^{-5} & -C_{1}\left(x_{1}+r\right)^{-2}\left|x^{\prime}\right|\left(r^{2}-\left|x^{\prime}\right|^{2}\right)^{-2} \\
& +\left(2 c\left|x^{\prime}\right|^{2}-C_{2} r^{4}\right)\left(r^{2}-\left|x^{\prime}\right|^{2}\right)^{-4}
\end{aligned}
$$

with constants $C_{1}, C_{2}$ independent of $r$. The negative term is bounded by $C_{3}\left(x_{1}+r\right)^{-4}+$ $c\left|x^{\prime}\right|^{2}\left(r^{2}-\left|x^{\prime}\right|^{2}\right)^{-4}$, hence

$$
\Delta_{\omega} v / v(x) \geqslant c\left(x_{1}+r\right)^{-5}+\left(c\left|x^{\prime}\right|^{2}-C_{2} r^{4}\right)\left(r^{2}-\left|x^{\prime}\right|^{2}\right)^{-4}
$$

The last term is negative only when $\left|x^{\prime}\right|<C_{4} r^{2}$, in which case it is bounded by $C_{5} r^{-4}<$ $c\left(x_{1}+r\right)^{-5}$. Hence $v$ is strongly $\omega$-subharmonic on $U$.

Next, assume that $U$ and $W$ are connected. Then $U \cup W$ is connected. Fix a point $a \in W \backslash \bar{U}$. If $z_{0} \in U$ is given, we choose a path $\Gamma \subset U \cup W$ from $z_{0}$ to $a$ which is piecewise linear with respect to holomorphic coordinate patches. Then we can find a finite sequence of cylinders $\left(U_{j}, W_{j}\right)$ of the type described above, $1 \leqslant j \leqslant N$, whose axes are segments contained in $\Gamma$, such that

$$
U_{j} \cup W_{j} \subset U \cup W, \quad \bar{W}_{j} \subset U_{j+1} \quad \text { and } \quad z_{0} \in U_{0}, \quad a \in W_{N} \subset W \backslash \bar{U}
$$

For each such pair, we have a function $v_{j} \in \mathscr{C}^{\infty}(M)$ with support in $\bar{U}_{j} \cup \bar{W}_{j}, v_{j} \geqslant 0$, strongly $\omega$-subharmonic and $>0$ on $U_{j}$. By induction, we can find constants $C_{j}>0$ such that $v_{0}+C_{1} v_{1}+\cdots+C_{j} v_{j}$ is strongly $\omega$-subharmonic on $U_{0} \cup \ldots \cup U_{j}$ and $\omega$-subharmonic on $M \backslash \bar{W}_{j}$. Then

$$
w_{z_{0}}=v_{0}+C_{1} v_{1}+\ldots+C_{N} v_{N} \geqslant 0
$$

is $\omega$-subharmonic on $U$ and strongly $\omega$-subharmonic $>0$ on a neighborhood $\Omega_{0}$ of the given point $z_{0}$. Select a denumerable covering of $U$ by such neighborhoods $\Omega_{p}$ and set $v(z)=\sum \varepsilon_{p} w_{z_{p}}(z)$ where $\varepsilon_{p}$ is a sequence converging sufficiently fast to 0 so that $v \in \mathscr{C}^{\infty}(M, \mathbb{R})$. Then $v$ has the required properties.

In the general case, we find for each pair $\left(U_{s}, W_{t(s)}\right)$ a function $v_{s}$ with support in $\bar{U}_{s} \cup \bar{W}_{t(s)}$, strongly $\omega$-subharmonic and $>0$ on $U_{s}$. Any convergent series $v=\sum \varepsilon_{s} v_{s}$ yields a function with the desired properties.
(3.4) Lemma. Let $X$ be a connected, locally connected and locally compact topological space. If $U$ is a relatively compact open subset of $X$, we let $\widetilde{U}$ be the union of $U$ with all compact connected components of $X \backslash U$. Then $\widetilde{U}$ is open and relatively compact in $X$, and $X \backslash \widetilde{U}$ has only finitely many connected components, all non compact.

Proof. A rather easy exercise of general topology. Intuitively, $\widetilde{U}$ is obtained by "filling the holes" of $U$ in $X$.
(3.5) Theorem ([Greene-Wu 1975]). Every n-dimensional connected non compact complex manifold $M$ has a strongly subharmonic exhaustion function with respect to any hermitian metric $\omega$. In particular, $M$ is strongly $n$-complete.

Proof. Let $\varphi \in \mathscr{C}^{\infty}(M, \mathbb{R})$ be an arbitrary exhaustion function. There exists a sequence of connected smoothly bounded open sets $\Omega_{\nu}^{\prime} \subset \subset M$ such that $\bar{\Omega}_{\nu}^{\prime} \subset \Omega_{\nu+1}^{\prime}$ and $M=\bigcup \Omega_{\nu}^{\prime}$. Let $\Omega_{\nu}=\widetilde{\Omega}_{\nu}^{\prime}$ be the relatively compact open set given by Lemma 3.4. Then $\bar{\Omega}_{\nu} \subset \Omega_{\nu+1}$, $M=\bigcup \Omega_{\nu}$ and $M \backslash \Omega_{\nu}$ has no compact connected component. We set

$$
U_{1}=\Omega_{2}, \quad U_{\nu}=\Omega_{\nu+1} \backslash \overline{\Omega_{\nu-2}} \quad \text { for } \nu \geqslant 2 .
$$

Then $\partial U_{\nu}=\partial \Omega_{\nu+1} \cup \partial \Omega_{\nu-2}$; any connected component $U_{\nu, s}$ of $U_{\nu}$ has its boundary $\partial U_{\nu, s} \not \subset \partial \Omega_{\nu-2}$, otherwise $\bar{U}_{\nu, s}$ would be open and closed in $M \backslash \Omega_{\nu-2}$, hence $\bar{U}_{\nu, s}$ would be a compact component of $M \backslash \Omega_{\nu-2}$. Therefore $\partial U_{\nu, s}$ intersects $\partial \Omega_{\nu+1} \subset$ $U_{\nu+1}$. If $\partial U_{\nu+1, t(s)}$ is a connected component of $U_{\nu+1}$ containing a point of $\partial U_{\nu, s}$, then $U_{\nu+1, t(s)} \cap U_{\nu, s} \neq \emptyset$ and $U_{\nu+1, t(s)} \backslash \bar{U}_{\nu, s} \neq \emptyset$. Lemma 7 implies that there is a nonnegative function $v_{\nu} \in \mathscr{C}^{\infty}(M, \mathbb{R})$ with support in $U_{\nu} \cup U_{\nu+1}$, which is strongly $\omega$-subharmonic on $U_{\nu}$. An induction yields constants $C_{\nu}$ such that

$$
\psi_{\nu}=\varphi+C_{1} v_{1}+\cdots+C_{\nu} v_{\nu}
$$

is strongly $\omega$-subharmonic on $\overline{\Omega_{\nu}} \subset U_{0} \cup \ldots \cup U_{\nu}$, thus $\psi=\varphi+\sum C_{\nu} v_{\nu}$ is a strongly $\omega$-subharmonic exhaustion function on $M$.

By an induction on the dimension, the above result can be generalized to an arbitrary complex space (or analytic scheme), as was first shown by T. Ohsawa.
(3.6) Theorem ([Ohsawa 1984]). Let $X$ be a complex space of maximal dimension $n$.
a) $X$ is always strongly $(n+1)$-complete.
b) If $X$ has no compact irreducible component of dimension $n$, then $X$ is strongly $n$ complete.
c) If $X$ has only finitely many irreducible components of dimension $n$, then $X$ is strongly $n$-convex.

Proof. We prove a) and b) by induction on $n=\operatorname{dim} X$. For $n=0$, property b) is void and a) is obvious (any function can then be considered as strongly 1-convex). Assume that a) has been proved in dimension $\leqslant n-1$. Let $X^{\prime}$ be the union of $X_{\text {sing }}$ and of the irreducible components of $X$ of dimension at most $n-1$, and $M=X \backslash X^{\prime}$ the $n$-dimensional part of $X_{\text {reg }}$. As $\operatorname{dim} X^{\prime} \leqslant n-1$, the induction hypothesis shows that $X^{\prime}$ is strongly $n$-complete. By Th. 2.13, there exists a strongly $n$-convex exhaustion function
$\varphi^{\prime}$ on a neighborhood $V^{\prime}$ of $X^{\prime}$. Take a closed neighborhood $\bar{V} \subset V^{\prime}$ and an arbitrary exhaustion $\varphi$ on $X$ that extends $\varphi^{\prime} \bar{V}$. Since every function on a $n$-dimensional manifold is strongly $(n+1)$-convex, we conclude that $X$ is at worst $(n+1)$-complete, as stated in a).

In case b), the hypothesis means that the connected components $M_{j}$ of $M=X \backslash X^{\prime}$ have non compact closure $\bar{M}_{j}$ in $X$. On the other hand, Lemma 3.1 shows that there exists a hermitian metric $\omega$ on $M$ such that $\varphi_{\mid M \cap V}$ is strongly $\omega$-subharmonic. Consider the open sets $U_{j, \nu} \subset M_{j}$ provided by Lemma 3.7 below. By the arguments already used in Th. 3.5, we can find a strongly $\omega$-subharmonic exhaustion $\psi=\varphi+\sum_{j, \nu} C_{j, \nu} v_{j, \nu}$ on $X$, with $v_{j, \nu}$ strongly $\omega$-subharmonic on $U_{j, \nu}, \operatorname{Supp} v_{j, \nu} \subset U_{j, \nu} \cup U_{j, \nu+1}$ and $C_{j, \nu}$ large. Then $\psi$ is strongly $n$-convex on $X$.
(3.7) Lemma. For each $j$, there exists a sequence of open sets $U_{j, \nu} \subset \subset M_{j}, \nu \in \mathbb{N}$, such that
a) $M_{j} \backslash V^{\prime} \subset \bigcup_{\nu} U_{j, \nu}$ and $\left(U_{j, \nu}\right)$ is locally finite in $\bar{M}_{j}$;
b) for every connected component $U_{j, \nu, s}$ of $U_{j, \nu}$ there is a connected component $U_{j, \nu+1, t(s)}$ of $U_{j, \nu+1}$ such that $U_{j, \nu+1, t(s)} \cap U_{j, \nu, s} \neq \emptyset$ and $U_{j, \nu+1, t(s)} \backslash \bar{U}_{j, \nu, s} \neq \emptyset$.

Proof. By Lemma 3.4 applied to the space $\bar{M}_{j}$, there exists a sequence of relatively compact connected open sets $\Omega_{j, \nu}$ in $\bar{M}_{j}$ such that $\bar{M}_{j} \backslash \Omega_{j, \nu}$ has no compact connected component, $\bar{\Omega}_{j, \nu} \subset \Omega_{j, \nu+1}$ and $\bar{M}_{j}=\bigcup \Omega_{j, \nu}$. We define a compact set $K_{j, \nu} \subset M_{j}$ and an open set $W_{j, \nu} \subset \bar{M}_{j}$ containing $K_{j, \nu}$ by

$$
K_{j, \nu}=\left(\bar{\Omega}_{j, \nu} \backslash \Omega_{j, \nu-1}\right) \backslash V^{\prime}, \quad W_{j, \nu}=\Omega_{j, \nu+1} \backslash \bar{\Omega}_{j, \nu-2}
$$

By induction on $\nu$, we construct an open set $U_{j, \nu} \subset \subset W_{j, \nu} \backslash X^{\prime} \subset M_{j}$ and a finite set $F_{j, \nu} \subset \partial U_{j, \nu} \backslash \bar{\Omega}_{j, \nu}$. We let $F_{j,-1}=\emptyset$. If these sets are already constructed for $\nu-1$, the compact set $K_{j, \nu} \cup F_{j, \nu-1}$ is contained in the open set $W_{j, \nu}$, thus contained in a finite union of connected components $W_{j, \nu, s}$. We can write $K_{j, \nu} \cup F_{j, \nu-1}=\bigcup L_{j, \nu, s}$ where $L_{j, \nu, s}$ is contained in $W_{j, \nu, s} \backslash X^{\prime} \subset M_{j}$. The open set $W_{j, \nu, s} \backslash X^{\prime}$ is connected and non contained in $\bar{\Omega}_{j, \nu} \cup L_{j, \nu, s}$, otherwise its closure $\bar{W}_{j, \nu, s}$ would have no boundary point $\in \partial \Omega_{j, \nu+1}$, thus would be open and compact in $\bar{M}_{j} \backslash \Omega_{j, \nu-2}$, contradiction. We select a point $a_{s} \in\left(W_{j, \nu, s} \backslash X^{\prime}\right) \backslash\left(\bar{\Omega}_{j, \nu} \cup L_{j, \nu, s}\right)$ and a smoothly bounded connected open set $U_{j, \nu, s} \subset \subset W_{j, \nu, s} \backslash X^{\prime}$ containing $L_{j, \nu, s}$ with $a_{s} \in \partial U_{j, \nu, s}$. Finally, we set $U_{j, \nu}=\bigcup_{s} U_{j, \nu, s}$ and let $F_{j, \nu}$ be the set of all points $a_{s}$. By construction, we have $U_{j, \nu} \supset K_{j, \nu} \cup F_{j, \nu-1}$, thus $\bigcup U_{j, \nu} \supset \bigcup K_{j, \nu}=M_{j} \backslash V^{\prime}$, and $\partial U_{j, \nu, s} \ni a_{s}$ with $a_{s} \in F_{j, \nu} \subset U_{j, \nu+1}$. Property b) follows.

Proof of Theorem 3.6 c) (end). Let $Y \subset X$ be the union of $X_{\text {sing }}$ with all irreducible components of $X$ that are non compact or of dimension $<n$. Then $\operatorname{dim} Y \leqslant n-1$, so $Y$ is $n$-convex and Th. 2.13 implies that there is an exhaustion function $\psi \in \mathscr{C}^{\infty}(X, \mathbb{R})$ such that $\psi$ is strongly $n$-convex on a neighborhood $V$ of $Y$. Then the complement $K=X \backslash V$ is compact and $\psi$ is strongly $n$-convex on $X \backslash K$.
(3.8) Proposition. Let $M$ be a connected non compact $n$-dimensional complex manifold and $U$ an open subset of $M$. Then $U$ is $n$-Runge complete in $M$ if and only if $M \backslash U$ has no compact connected component.

Proof. First observe that a strongly $n$-convex function cannot have any local maximum, so it satisfies the maximum principle. If $M \backslash U$ has a compact connected component $T$, then $T$ has a compact neighborhood $L$ in $M$ such that $\partial L \subset U$. We have $\max _{L} \psi=\max _{\partial L} \psi$ for every strongly $n$-convex function, thus $\partial L \subset M_{b}$ implies $L \subset M_{b}$; thus we cannot find a sublevel set $M_{b}$ such that $\partial L \subset M_{b} \subset \subset U$, and $U$ is not $n$-Runge in $M$.

On the other hand, assume that $M \backslash U$ has no compact connected component and let $L$ be a compact subset of $U$. Let $\omega$ be any hermitian metric on $M$ and $\varphi$ a strongly $\omega$-subharmonic exhaustion function on $M$. Set $b=1+\sup _{L} \varphi$ and

$$
P=\{x \in M \backslash U ; \varphi(x) \leqslant b\} .
$$

As $M \backslash U$ has no compact connected component, all its components $T_{\alpha}$ contain a point $y_{\alpha}$ in

$$
W=\{x \in X ; \varphi(x)>b+1\} .
$$

For every point $x \in P$ with $x \in T_{\alpha}$, there exists a connected open set $V_{x} \subset \subset M \backslash L$ containing $x$ such that $\partial V_{x} \ni y_{\alpha}(M \backslash L$ is a neighborhood of $M \backslash U$ and we can consider a tubular neighborhood of a path from $x$ to $y_{\alpha}$ in $\left.M \backslash L\right)$. The compact set $P$ can be covered by a finite number of open sets $V_{x_{j}}$. Then Lemma 3.3 yields functions $v_{j}$ with support in $\bar{V}_{x_{j}} \cup \bar{W}$ which are strongly $\omega$-subharmonic on $V_{x_{j}}$. Let $\chi$ be a convex increasing function such that $\chi(t)=0$ on $]-\infty, b]$ and $\chi^{\prime}(t)>0$ on $] b,+\infty[$. Consider the function

$$
\psi=\varphi+\sum C_{j} v_{j}+\chi \circ \varphi .
$$

First, choose $C_{j}$ large enough so that $\psi \geqslant b$ on $P$. Then choose $\chi$ increasing fast enough so that $\psi$ is strongly $\omega$-subharmonic on $\bar{W}$. Then $\psi$ is a strongly $n$-convex exhaustion function on $M$, and as $\psi \geqslant \varphi$ on $M$ and $\psi=\varphi$ on $L$, we see that

$$
L \subset\{x \in M ; \psi(x)<b\} \subset U
$$

This proves that $U$ is $n$-Runge complete in $M$.

## §4. Andreotti-Grauert Finiteness Theorems

## $\S 4$.A. Case of Vector Bundles over Manifolds

The crucial point in the proof of the Andreotti-Grauert theorems is the following special case, which is easily obtained by the methods of chapter 8 .
(4.1) Proposition. Let $M$ be a strongly $q$-complete manifold with $q \geqslant 1$, and $E$ a holomorphic vector bundle over M. Then:
a) $H^{k}(M, \mathscr{O}(E))=0$ for $k \geqslant q$.
b) Let $U$ be a q-Runge complete open subset of $M$. Every $d^{\prime \prime}$-closed form $h \in$ $\mathscr{C}_{0, q-1}^{\infty}(U, E)$ can be approximated uniformly with all derivatives on every compact subset of $U$ by a sequence of global $d^{\prime \prime}$-closed forms $\widetilde{h}_{\nu} \in \mathscr{C}_{0, q-1}^{\infty}(M, E)$.

Proof. We replace $E$ by $\widetilde{E}=\Lambda^{n} T M \otimes E$; then we can work with forms of bidegree $(n, k)$ instead of $(0, k)$. Let $\psi$ be a strongly $q$-convex exhaustion function on $M$ and $\omega$
the metric given by Lemma 3.1. Select a function $\rho \in \mathscr{C}^{\infty}(M, \mathbb{R})$ which increases rapidly at infinity so that the hermitian metric $\widetilde{\omega}=e^{\rho} \omega$ is complete on $M$. Denote by $E_{\chi}$ the bundle $E$ endowed with the hermitian metric obtained by multiplication of a fixed metric of $E$ by the weight $\exp (-\rho \circ \psi)$ where $\chi \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$ is a convex increasing function. We apply Th. VIII- 4.5 for the bundle $E_{\chi}$ over the complete hermitian manifold ( $M, \widetilde{\omega}$ ). Then

$$
i c\left(E_{\chi}\right)=i c(E)+i d^{\prime} d^{\prime \prime}(\chi \circ \psi) \otimes \operatorname{Id}_{E} \geqslant_{\mathrm{Nak}} i c(E)+\chi^{\prime} \circ \psi i d^{\prime} d^{\prime \prime} \psi \otimes \operatorname{Id}_{E}
$$

The eigenvalues of $i d^{\prime} d^{\prime \prime} \psi$ with respect to $\widetilde{\omega}$ are $e^{-\rho} \gamma_{j}$, so Lemma VII-7.2 and Prop. VI-8.3 yield

$$
\begin{aligned}
{\left[i c\left(E_{\chi}\right), \Lambda\right]+T_{\widetilde{\omega}} } & \geqslant[i c(E), \Lambda]+T_{\widetilde{\omega}}+\chi^{\prime} \circ \psi\left[i d^{\prime} d^{\prime \prime} \psi, \Lambda\right] \otimes \operatorname{Id}_{E} \\
& \geqslant[i c(E), \Lambda]+T_{\widetilde{\omega}}+\chi^{\prime} \circ \psi e^{-\rho}\left(\gamma_{1}+\cdots+\gamma_{k}\right) \otimes \operatorname{Id}_{E}
\end{aligned}
$$

when this curvature tensor acts on $(n, k)$-forms. For $k \geqslant q$, we have

$$
\gamma_{1}+\cdots+\gamma_{k} \geqslant 1-(q-1) \varepsilon>0 \quad \text { if } \varepsilon \leqslant 1 / q
$$

We choose $\chi_{0}$ increasing fast enough so that all the eigenvalues of the above curvature tensor are $\geqslant 1$ when $\chi=\chi_{0}$. Then for every $g \in \mathscr{C}_{n, k}^{\infty}(M, E)$ with $D^{\prime \prime} g=0$ the equation $D^{\prime \prime} f=g$ can be solved with an estimate

$$
\int_{M}|f|^{2} e^{-\chi \circ \psi} d V \leqslant \int_{M}|g|^{2} e^{-\chi \circ \psi} d V
$$

where $\chi=\chi_{0}+\chi_{1}$ and where $\chi_{1}$ is a convex increasing function chosen so that the integral of $g$ converges. This gives a). In order to prove b), let $h \in \mathscr{C}_{n, q-1}^{\infty}(U, E)$ be such that $D^{\prime \prime} h=0$ and let $L$ be an arbitrary compact subset of $U$. Thanks to Def. 2.14, we can choose $\psi$ such that there is a sublevel set $M_{b}$ with $L \subset M_{b} \subset \subset U$. Select $b_{0}<b$ so that $L \subset M_{b_{0}}$, and let $\theta \in \mathscr{C}^{\infty}(\mathbb{R}, \mathbb{R})$ be a convex increasing function such that $\theta=0$ on $]-\infty, b_{0}[$ and $\theta \geqslant 1$ on $] b,+\infty[$. Let $\eta \in \mathscr{D}(U)$ be a cut-off function such that $\eta=1$ on $M_{b}$. We solve the equation $D^{\prime \prime} f=g$ for $g=D^{\prime \prime}(\eta h)$ with the weight $\chi=\chi_{0}+\nu \theta \circ \psi$ and let $\nu$ tend to infinity. As $g$ has compact support in $U \backslash M_{b}$ and $\chi \circ \psi \geqslant \chi_{0} \circ \psi+\nu$ on this set, we find a solution $f_{\nu}$ such that

$$
\int_{M_{b_{0}}}\left|f_{\nu}\right|^{2} e^{-\chi_{0} \circ \psi} d V \leqslant \int_{M}\left|f_{\nu}\right|^{2} e^{-\chi \circ \psi} d V \leqslant \int_{U \backslash M_{b}}|g|^{2} e^{-\chi \circ \psi} d V \leqslant C e^{-\nu}
$$

thus $f_{\nu}$ converges to 0 in $L^{2}\left(M_{b_{0}}\right)$ and $h_{\nu}=\eta h-f_{\nu} \in \mathscr{C}_{n, q-1}^{\infty}(M, E)$ is a $D^{\prime \prime}$-closed form converging to $h$ in $L^{2}\left(M_{b_{0}}\right)$. However, if we choose the minimal solution such that $\delta_{\chi}^{\prime \prime} f_{\nu}=0$ as in Rem. VIII-4.6, we get $\Delta_{\chi}^{\prime \prime} f_{\nu}=\delta_{\chi}^{\prime \prime} g$ on $M$ and in particular $\Delta_{\chi_{0}}^{\prime \prime} f_{\nu}=0$ on $M_{b_{0}}$. Gårding's inequality VI-3.3 applied to the elliptic operator $\Delta_{\chi_{0}}^{\prime \prime}$ shows that $f_{\nu}$ converges to 0 with all derivatives on $L$, hence $h_{\nu}$ converges to $h$ on $L$. Now, replace $L$ by an exhaustion $L_{\nu}$ of $U$ by compact sets; some diagonal subsequence $h_{\nu}$ converges to $h$ in $\mathscr{C}_{n, q-1}^{\infty}(U, E)$.

## §4.B. A Local Vanishing Result for Sheaves

Let $\left(X, \mathscr{O}_{X}\right)$ be an analytic scheme and $\mathscr{T}$ a coherent sheaf of $\mathscr{O}_{X}$-modules. We wish to extend Prop. 4.1 to the cohomology groups $H^{k}(X, \mathscr{P})$. The first step is to show that the result holds on small open sets, and this is done by means of local resolutions of $\mathscr{S}$.

For a given point $x \in X$, we choose a patch $\left(A, \mathscr{O}_{\Omega} / \mathscr{F}\right)$ of $X$ containing $x$, where $A$ is an analytic subset of $\Omega \subset \mathbb{C}^{N}$ and $\mathscr{g}$ a sheaf of ideals with zero set $A$. Let $i_{A}: A \longrightarrow \Omega$ be the inclusion. Then $\left(i_{A}\right)_{\star} \mathscr{S}$ is a coherent $\mathscr{O}_{\Omega}$-module supported on $A$. In particular there is a neighborhood $W_{0} \subset \Omega$ of $x$ and a surjective sheaf morphism

$$
\mathscr{O}^{p_{0}} \longrightarrow\left(i_{A}\right)_{\star} \mathscr{S} \quad \text { on } \quad W_{0}, \quad\left(u_{1}, \ldots, u_{p_{0}}\right) \longmapsto \sum_{1 \leqslant j \leqslant p_{0}} u_{j} G_{j}
$$

where $G_{1}, \ldots, G_{p_{0}} \in \mathscr{S}\left(A \cap W_{0}\right)$ are generators of $\left(i_{A}\right)_{\star} \mathscr{S}$ on $W_{0}$. If we repeat the procedure inductively for the kernel of the above surjective morphism, we get a homological free resolution of $\left(i_{A}\right)_{\star} \mathscr{S}$ :

$$
\begin{equation*}
\mathscr{O}^{p_{l}} \longrightarrow \cdots \longrightarrow \mathscr{O}^{p_{1}} \longrightarrow \mathscr{O}^{p_{0}} \longrightarrow\left(i_{A}\right)_{\star} \mathscr{S} \longrightarrow 0 \quad \text { on } W_{l} \tag{4.3}
\end{equation*}
$$

of arbitrary large length $l$, on neighborhoods $W_{l} \subset W_{l-1} \subset \ldots \subset W_{0}$. In particular, after replacing $\Omega$ by $W_{2 N}$ and $A$ by $A \cap W_{2 N}$, we may assume that $\left(i_{A}\right)_{\star} \mathscr{S}$ has a resolution of length $2 N$ on $\Omega$. In this case, we shall say that $A \subset \Omega$ is a $\mathscr{T}$-distinguished patch of $X$.
(4.4) Lemma. Let $A \subset \Omega$ be a $\mathscr{S}$-distinguished patch of $X$ and $U$ a strongly $q$-convex open subset of $A$. Then

$$
H^{k}(U, \mathscr{P})=0 \quad \text { for } \quad k \geqslant q .
$$

Proof. Theorem 2.13 shows that there exists a strongly $q$-convex open set $V \subset \Omega$ such that $U=A \cap V$. Let us denote by $\mathscr{E}^{l}$ the kernel of $\mathscr{O}^{p_{l}} \longrightarrow \mathbb{O}^{p_{l-1}}$ for $l \geqslant 1$ and $\mathscr{E}^{0}=\operatorname{ker}\left(\mathscr{C}^{p_{0}} \longrightarrow\left(i_{A}\right)_{\star} \mathscr{Y}\right)$. There are exact sequences

$$
\begin{aligned}
& 0 \longrightarrow \mathscr{L}^{0} \longrightarrow \mathscr{O}^{p_{0}} \longrightarrow\left(i_{A}\right)_{\star} \mathscr{Y} \longrightarrow 0 \\
& 0 \longrightarrow \mathscr{L}^{l} \longrightarrow \mathscr{O}^{p_{l}} \longrightarrow \mathscr{L}^{l-1} \longrightarrow 0, \quad 1 \leqslant l \leqslant 2 N
\end{aligned}
$$

For $k \geqslant q$, Prop. 4.1 a) gives $H^{k}\left(V, \mathscr{O}^{p_{l}}\right)=0$, therefore we get

$$
H^{k}(U, \mathscr{Y}) \simeq H^{k}\left(V,\left(i_{A}\right)_{\star} \mathscr{Y}\right) \simeq H^{k+1}\left(V, \mathscr{Z}^{0}\right) \simeq \ldots \simeq H^{k+2 N+1}\left(V, \mathscr{Z}^{2 N}\right)
$$

and the last group vanishes because topdim $V \leqslant \operatorname{dim}_{\mathbb{R}} V=2 N$.

## $\S$ 4.C. Topological Structure on Spaces of Sections and on Cohomology Groups

Let $V \subset \Omega$ be a strongly 1-complete open set relatively to a $\mathscr{S}$-distinguished patch $A \subset \Omega$ and let $U=A \cap V$. By the proof of Lemma 4.4, we have

$$
H^{1}\left(V, \mathscr{L}^{0}\right) \simeq H^{2 N+1}\left(V, \mathscr{L}^{2 N}\right)=0
$$

hence we get an exact sequence

$$
\begin{equation*}
0 \longrightarrow \mathscr{L}^{0}(V) \longrightarrow \mathscr{O}^{p_{0}}(V) \longrightarrow \mathscr{S}(U) \longrightarrow 0 \tag{4.5}
\end{equation*}
$$

We are going to show that the Fréchet space structure on $\mathbb{O}^{p_{0}}(V)$ induces a natural Fréchet space structure on the groups of sections of $\mathscr{S}$ over any open subset. We first note that $\mathscr{E}^{0}(V)$ is closed in $\mathscr{C}^{p_{0}}(V)$. Indeed, let $f_{\nu} \in \mathscr{Z}^{0}(V)$ be a sequence converging to a limit $f \in \mathscr{O}^{p_{0}}(V)$ uniformly on compact subsets of $V$. For every $x \in V$, the germs $\left(f_{\nu}\right)_{x}$ converge to $f_{x}$ with respect to the topology defined by (1.4) on $\mathscr{O}^{p_{0}}$. As $\mathscr{E}_{x}^{0}$ is closed in $\mathscr{O}_{x}^{p_{0}}$ in view of Th. 1.5 b ), we get $f_{x} \in \mathscr{L}_{x}^{0}$ for all $x \in V$, thus $f \in \mathscr{L}^{0}(V)$.
(4.6) Proposition. The quotient topology on $\mathscr{S}(U)$ is independent of the choices made above.

Proof. For a smaller set $U^{\prime}=A \cap V^{\prime}$ where $V^{\prime}$ is a strongly 1-convex open subset of $V$, the restriction map $\mathscr{O}^{p_{0}}(V) \longrightarrow \mathscr{O}^{p_{0}}\left(V^{\prime}\right)$ is continuous, thus $\mathscr{S}(U) \longrightarrow \mathscr{I}\left(U^{\prime}\right)$ is continuous. If $\left(V_{\alpha}\right)$ is a countable covering of $V$ by such sets and $U_{\alpha}=A \cap V_{\alpha}$, we get an injection of $\mathscr{P}(U)$ onto the closed subspace of the product $\prod \mathscr{P}\left(U_{\alpha}\right)$ consisting of families which are compatible in the intersections. Therefore, the Fréchet topology induced by the product coincides with the original topology of $\mathscr{S}(U)$. If we choose other generators $H_{1}, \ldots, H_{q_{0}}$ for $\left(i_{A}\right)_{\star} \mathscr{P}$, the germs $H_{j, x}$ can be expressed in terms of the $G_{j, x}$ 's, thus we get a commutative diagram

provided that $U$ and $V$ are small enough. If we express the generators $G_{j}$ in terms of the $H_{j}$ 's, we find a similar diagram with opposite vertical arrows and we conclude easily that the topology obtained in both cases is the same. Finally, it remains to show that the topology of $\mathscr{S}(U)$ is independent of the embedding $A \subset \Omega$ near a given point $x \in X$. We compare the given embedding with the Zariski embedding $(A, x) \subset \Omega^{\prime}$ of minimal dimension $d$. After shrinking $A$ and changing coordinates, we may assume $\Omega=\Omega^{\prime} \times \mathbb{C}^{N-d}$ and that the embedding $i_{A}: A \longrightarrow \Omega$ is the composite of $i_{A}^{\prime}: A \longrightarrow \Omega^{\prime}$ and of the inclusion $j: \Omega^{\prime} \longrightarrow \Omega^{\prime} \times\{0\} \subset \Omega$. For $V^{\prime} \subset \Omega^{\prime}$ sufficient small and $U^{\prime}=A \cap V^{\prime}$, we have a surjective map $G^{\prime}: \mathscr{O}^{p_{0}}\left(V^{\prime}\right) \longrightarrow \mathscr{S}\left(U^{\prime}\right)$ obtained by choosing generators $G_{j}^{\prime}$ of $\left(i_{A}^{\prime}\right)^{\star} \mathscr{S}$ on a neighborhood of $x$ in $\Omega^{\prime}$. Then we consider the open set $V=V^{\prime} \times \mathbb{C}^{N-d} \subset \Omega$ and the surjective map onto $\mathscr{P}\left(U^{\prime}\right)$ equal to the composite

$$
\mathscr{O}^{p_{0}}(V) \xrightarrow{j^{\star}} \mathscr{O}^{p_{0}}\left(V^{\prime}\right) \xrightarrow{G^{\prime}} \mathscr{S}(U)
$$

This map corresponds to a choice of generators $G_{j} \in\left(i_{A}\right)^{\star} \mathscr{S}(V)$ equal to the functions $G_{j}^{\prime}$, considered as functions independent of the last variables $z_{d+1}, \ldots, z_{N}$. Since $j^{\star}$ is open, it is obvious that the quotient topology on $\mathscr{S}\left(U^{\prime}\right)$ is the same for both embeddings.

Now, there is a natural topology on the cohomology groups $H^{k}(X, \mathscr{Y})$. In fact, let ( $U_{\alpha}$ ) be a countable covering of $X$ by strongly 1-complete open sets, such that each $U_{\alpha}$ is contained in a $\mathscr{T}$-distinguished patch. Since the intersections $U_{\alpha_{0} \ldots \alpha_{k}}$ are again strongly 1-complete, the covering $\mathscr{U}$ is acyclic by Lemma 4.4 and Leray's theorem shows that $H^{k}(X, \mathscr{Y})$ is isomorphic to $\check{H}^{q}(\mathscr{U}, \mathscr{Y})$. We consider the product topology on the spaces of Čech cochains $C^{k}(\mathscr{U}, \mathscr{Y})=\prod \mathscr{P}\left(U_{\alpha_{0} \ldots \alpha_{k}}\right)$ and the quotient topology on $\check{H}^{k}(\mathscr{U}, \mathscr{S})$. It is clear that $\check{H}^{0}(\mathscr{U}, \mathscr{S})$ is a Fréchet space; however the higher cohomology groups $\check{H}^{k}(\mathscr{U}, \mathscr{\mathscr { S }})$ need not be Hausdorff because the coboundary groups may be non closed in the cocycle
groups. The resulting topology on $H^{k}(X, \mathscr{Y})$ is independent of the choice of the covering: in fact we only have to check that the bijective continuous map $\check{H}^{k}(\mathscr{U}, \mathscr{S}) \longrightarrow \check{H}^{k}\left(\mathscr{U}^{\prime}, \mathscr{S}\right)$ is a topological isomorphism if $\mathscr{U}^{\prime}$ is a refinement of $\mathscr{U}$, and this follows from Rem. 1.10 applied to the morphism of Čech complexes $C^{\bullet}(\mathscr{U}, \mathscr{S}) \longrightarrow C^{\bullet}\left(\mathscr{U}^{\prime}, \mathscr{S}\right)$.

Finally, observe that when $\mathscr{S}$ is the locally free sheaf associated to a holomorphic vector bundle $E$ on a smooth manifold $X$, the topology on $H^{k}(X, \mathscr{O}(E))$ is the same as the topology associated to the Fréchet space structure on the Dolbeault complex $\left(C_{0, \bullet}^{\infty}(X, E), d^{\prime \prime}\right)$ : by the analogue of formula (IV-6.11) we have a bijective continuous map

$$
\begin{aligned}
\check{H}^{k}(\vartheta, \overparen{U},(E)) & \longrightarrow H^{k}\left(\mathscr{C}_{0, \bullet}^{\infty}(X, E)\right) \\
\left\{\left(c_{\alpha_{0} \ldots \alpha_{k}}\right)\right\} & \longmapsto f(z)=\sum_{\alpha_{0}, \ldots, \alpha_{q}} c_{\alpha_{0} \ldots \alpha_{q}}(z) \theta_{\alpha_{q}} d^{\prime \prime} \theta_{\alpha_{0}} \wedge \ldots \wedge d^{\prime \prime} \theta_{\alpha_{q-1}}
\end{aligned}
$$

where $\left(\theta_{\alpha}\right)$ is a partition of unity subordinate to $\mathscr{U}$. As in Rem. 1.10, the continuity of the inverse follows by the open mapping theorem applied to the surjective map

$$
Z^{k}\left(C^{\bullet}(\mathscr{U}, \mathscr{O}(E))\right) \oplus \mathscr{C}_{0, k-1}^{\infty}(X, E) \longrightarrow Z^{k}\left(\mathscr{C}_{0, \bullet}^{\infty}(X, E)\right)
$$

We shall need a few simple additional results.
(4.7) Proposition. The following properties hold:
a) For every $x \in X$, the map $\mathscr{I}(X) \longrightarrow \mathscr{S}_{x}$ is continuous with respect to the topology of $\mathscr{S}_{x}$ defined by (1.4).
b) If $\mathscr{G}^{\prime}$ is a coherent analytic subsheaf of $\mathscr{T}$, the space of global sections $\mathscr{S}^{\prime}(X)$ is closed in $\mathscr{S}(X)$.
c) If $U^{\prime} \subset U$ are open in $X$, the restriction maps $H^{k}(U, \mathscr{Y}) \longrightarrow H^{k}\left(U^{\prime}, \mathscr{Y}\right)$ are continuous.
d) If $U^{\prime}$ is relatively compact in $U$, the restriction operator $\mathscr{P}(U) \longrightarrow \mathscr{I}\left(U^{\prime}\right)$ is compact.
e) Let $\mathscr{S} \longrightarrow \mathscr{S}^{\prime}$ be a morphism of coherent sheaves over $X$. Then the induced maps $H^{k}(X, \mathscr{P}) \longrightarrow H^{k}\left(X, \mathscr{S}^{\prime}\right)$ are continuous.

Proof. a) Let $V \subset \Omega$ be a strongly 1-convex open neighborhood of $x$ relatively to a $\mathscr{S}$-distinguished patch $A \subset \Omega$. The map $\mathscr{C}^{p_{0}}(V) \longrightarrow \mathscr{O}_{x}^{p_{0}}$ is continuous, and the same is true for $\mathscr{O}_{x}^{p_{0}} \longrightarrow \mathscr{S}_{x}$ by $\S 1$. Therefore the composite $\mathscr{O}^{p_{0}}(V) \longrightarrow \mathscr{S}_{x}$ and its factorization $\mathscr{S}(U) \longrightarrow \mathscr{S}_{x}$ are continuous.
b) is a consequence of the above property a) and of the fact that each stalk $\mathscr{S}_{x}^{\prime}$ is closed in $\mathscr{S}_{x}($ cf. 1.5 b$)$ ).
c) The restriction map $\mathscr{S}(U) \longrightarrow \mathscr{S}\left(U^{\prime}\right)$ is continuous, and the case of higher cohomology groups follows immediately.
d) Assume first that $U=A \cap V$ and $U^{\prime}=A \cap V^{\prime}$, where $A \subset \Omega$ is a $\mathscr{S}$-distinguished patch and $V^{\prime} \subset \subset V$ are strongly 1 -convex open subsets of $\Omega$. The operator $\mathscr{0}^{p_{0}}(V) \longrightarrow \mathscr{0}^{p_{0}}\left(V^{\prime}\right)$ is compact by Montel's theorem, thus $\mathscr{S}(U) \longrightarrow \mathscr{S}\left(U^{\prime}\right)$ is also compact. In the general
case, select a finite family of strongly 1-convex sets $U_{\alpha}^{\prime} \subset \subset U_{\alpha} \subset U$ such that $\left(U_{\alpha}^{\prime}\right)$ covers $\bar{U}^{\prime}$ and $U_{\alpha}$ is contained in some distinguished patch. There is a commutative diagram

where the right vertical arrow is a monomorphism and where the first arrow in the bottom line is compact. Thus $\mathscr{I}(U) \longrightarrow \mathscr{I}\left(U^{\prime}\right)$ is compact.
e) It is enough to check that $\mathscr{S}(U) \longrightarrow \mathscr{S}^{\prime}(U)$ is continuous, and for this we may assume that $U=A \cap V$ where $V$ is a small neighborhood of a given point $x$. Let $G_{1}, \ldots, G_{p_{0}}$ be generators of $\mathscr{S}_{x}, G_{1}^{\prime}, \ldots, G_{p_{0}}^{\prime}$ their images in $\mathscr{S}_{x}^{\prime}$. Complete these elements in order to obtain a system of generators $\left(G_{1}^{\prime}, \ldots, G_{q_{0}}^{\prime}\right)$ of $\mathscr{S}_{x}^{\prime}$. For $V$ small enough, the map $\mathscr{S}(U) \longrightarrow \mathscr{S}^{\prime}(U)$ is induced by the inclusion $\mathscr{\mathscr { O }}^{p_{0}}(V) \longrightarrow \mathscr{O}^{p_{0}}(V) \times\{0\} \subset \mathscr{O}^{q_{0}}(V)$, hence continuous.

## § 4.D. Cartan-Serre Finiteness Theorem

The above results enable us to prove a finiteness theorem for cohomology groups over compact analytic schemes.
(4.8) Theorem (Cartan-Serre). Let $\mathscr{S}$ be a coherent analytic sheaf over an analytic scheme $\left(X, \mathfrak{Q}_{X}\right)$. If $X$ is compact, all cohomology groups $H^{k}(X, \mathscr{S})$ are finite dimensional (and Hausdorff ).

Proof. There exist finitely many strongly 1-complete open sets $U_{\alpha}^{\prime} \subset \subset U_{\alpha}$ such that each $U_{\alpha}$ is contained in some $\mathscr{\mathscr { S }}$-distinguished patch and such that $\bigcup U_{\alpha}^{\prime}=X$. By Prop. 4.7 d ), the restriction map on Čech cochains

$$
C^{\bullet}(\mathscr{U}, \mathscr{S}) \longrightarrow C^{\bullet}\left(\mathscr{U}^{\prime}, \mathscr{Y}\right)
$$

defines a compact morphism of complexes of Fréchet spaces. As the coverings $\mathscr{U}=\left(U_{\alpha}\right)$ and $U^{\prime}=\left(U_{\alpha}^{\prime}\right)$ are acyclic by 4.4 , the induced map

$$
\check{H}^{k}(\mathscr{U}, \mathscr{S}) \longrightarrow \check{H}^{k}\left(\mathscr{U}^{\prime}, \mathscr{S}\right)
$$

is an isomorphism, both spaces being isomorphic to $H^{k}(X, \mathscr{Y})$. We conclude by Schwartz' theorem 1.9.

## §4.E. Local Approximation Theorem

We show that a local analogue of the approximation result 4.1 b ) holds for a sheaf $\mathscr{S}$ over an analytic scheme ( $X, \mathscr{O}_{X}$ ).
(4.9) Lemma. Let $A \subset \Omega$ be a $\mathscr{G}$-distinguished patch of $X$, and $U^{\prime} \subset U \subset A$ open subsets such that $U^{\prime}$ is $q$-Runge complete in $U$. Then the restriction map

$$
H^{q-1}(U, \mathscr{Y}) \longrightarrow H^{q-1}\left(U^{\prime}, \mathscr{S}\right)
$$

has a dense range.

Proof. Let $L$ be an arbitrary compact subset of $U^{\prime}$. Proposition 2.16 applied with $Y=U$ embedded in some neighborhood in $\Omega$ shows that there is a neighborhood $V$ of $U$ in $\Omega$ such that $A \cap V=U$ and a strongly $q$-convex function $\psi$ on $V$ such that $L \subset U_{b} \subset \subset U^{\prime}$ for some $U_{b}=A \cap V_{b}$. The proof of Lemma 4.4 gives $H^{q}\left(V, \mathscr{L}^{0}\right)=H^{q}\left(V_{b}, \mathscr{E}^{0}\right)=0$ and the cohomology exact sequences of $0 \rightarrow \mathscr{Z}^{0} \rightarrow \mathscr{O}^{p_{0}} \rightarrow i_{A}^{\star} \mathscr{Y} \rightarrow 0$ over $V$ and $V_{b}$ yield a commutative diagram of continuous maps

where the horizontal arrows are surjective. Since $V_{b}$ is $q$-Runge complete in $V$, the left vertical arrow has a dense range by Prop. 4.1 b ). As $U^{\prime}$ is the union of an increasing sequence of sets $U_{b_{\nu}}$, we only have to show that the range remains dense in the inverse limit $H^{q-1}\left(U^{\prime}, \mathscr{P}\right)$. For that, we apply Property 1.11 d$)$ on a suitable covering of $U$. Let ${ }^{W} W$ be a countable basis of the topology of $U$, consisting of strongly 1-convex open subsets contained in $\mathscr{G}$-distinguished patches. We let ${ }^{\prime} W^{\prime}$ (resp. ${ }^{W} W_{\nu}$ ) be the subfamily of $W \in{ }^{W} W$ such that $W \subset \subset U^{\prime}$ (resp. $W \subset \subset U_{b_{\nu}}$ ). Then ${ }^{\mathscr{W}} \mathbb{W}^{\prime} W^{\prime}, W_{\nu}$ are acyclic coverings of $U, U^{\prime}, U_{b_{\nu}}$ and each restriction map $C^{\bullet}\left({ }^{\bullet} W, \mathscr{Y}\right) \longrightarrow C^{\bullet}\left({ }^{`} W_{\nu}, \mathscr{S}\right)$ is surjective. Property $1.11 \mathrm{~d})$ can thus be applied and the lemma follows.

## §4.F. Statement and Proof of the Andreotti-Grauert Theorem

(4.10) Theorem ([Andreotti-Grauert 1962]). Let $\mathscr{T}$ be a coherent analytic sheaf over a strongly $q$-convex analytic scheme $\left(X, \mathscr{O}_{X}\right)$. Then
a) $H^{k}(X, \mathscr{P})$ is Hausdorff and finite dimensional for $k \geqslant q$.

Moreover, let $U$ be a $q$-Runge open subset of $X, q \geqslant 1$. Then
b) the restriction map $H^{k}(X, \mathscr{P}) \rightarrow H^{k}(U, \mathscr{S})$ is an isomorphism for $k \geqslant q$;
c) the restriction map $H^{q-1}(X, \mathscr{P}) \rightarrow H^{q-1}(U, \mathscr{S})$ has a dense range.

The compact case $q=0$ of 4.10 a) is precisely the Cartan-Serre finiteness theorem. For $q \geqslant 1$, the special case when $X$ is strongly $q$-complete and $U=\emptyset$ yields the following very important consequence.
(4.11) Corollary. If $X$ is strongly $q$-complete, then

$$
H^{k}(X, \mathscr{S})=0 \quad \text { for } \quad k \geqslant q .
$$

Assume that $q \geqslant 1$ and let $\psi$ be a smooth exhaustion on $X$ that is strongly $q$-convex on $X \backslash K$. We first consider sublevel sets $X_{d} \supset X_{c} \supset K, d>c$, and verify assertions $4.10 \mathrm{~b}), \mathrm{c})$ for all restriction maps

$$
H^{k}\left(X_{d}, \mathscr{Y}\right) \longrightarrow H^{k}\left(X_{c}, \mathscr{S}\right), \quad k \geqslant q-1 .
$$

The basic idea, already contained in [Andreotti-Grauert 1962], is to deform $X_{c}$ into $X_{d}$ through a sequence of strongly $q$-convex open sets $\left(G_{j}\right)$ such that $G_{j+1}$ is obtained from $G_{j}$ by making a small bump.
(4.12) Lemma. There exist a sequence of strongly $q$-convex open sets $G_{0} \subset \ldots \subset G_{s}$ and a sequence of strongly $q$-complete open sets $U_{0}, \ldots, U_{s-1}$ in $X$ such that
a) $G_{0}=X_{c}, \quad G_{s}=X_{d}, \quad G_{j+1}=G_{j} \cup U_{j} \quad$ for $0 \leqslant j \leqslant s-1$;
b) $G_{j}=\left\{x \in X ; \psi_{j}(x)<c_{j}\right\}$ where $\psi_{j}$ is an exhaustion function on $X$ that is strongly $q$-convex on $X \backslash K$;
c) $U_{j}$ is contained in a $\mathscr{G}$-distinguished patch $A_{j} \subset \Omega_{j}$ of $X$;
d) $G_{j} \cap U_{j}$ is strongly $q$-complete and $q$-Runge complete in $U_{j}$.

Proof. There exists a finite covering of the compact set $\bar{X}_{d} \backslash X_{c}$ by $\mathscr{\mathscr { L }}$-distinguished patches $A_{j} \subset \Omega_{j}, 0 \leqslant j<s$, where $\Omega_{j} \subset \mathbb{C}^{N_{j}}$ is a euclidean ball and $K \cap A_{j}=\emptyset$. Let $\theta_{j} \in \mathscr{D}(X)$ be a family of functions such that $\operatorname{Supp} \theta_{j} \subset A_{j}, \theta_{j} \geqslant 0, \sum \theta_{j} \leqslant 1$ and $\sum \theta_{j}=1$ on a neighborhood of $\bar{X}_{d} \backslash X_{c}$. We can find $\varepsilon_{0}>0$ so small that

$$
\psi_{j}=\psi-\varepsilon \sum_{0 \leqslant k<j} \theta_{k}
$$

is still strongly $q$-convex on $X \backslash K$ for $0 \leqslant j \leqslant s$ and $\varepsilon \leqslant \varepsilon_{0}$. We have $\psi_{0}=\psi$ and $\psi_{s}=\psi-\varepsilon$ on $\bar{X}_{d} \backslash X_{c}$, thus

$$
G_{j}=\left\{x \in X ; \psi_{j}(x)<c\right\}, \quad 0 \leqslant j \leqslant s
$$

is an increasing sequence of strongly $q$-convex open sets such that $G_{0}=X_{c}, G_{s}=X_{c+\varepsilon}$. Moreover, as $\psi_{j+1}-\psi_{j}=-\varepsilon \theta_{j}$ has support in $A_{j}$, we have

$$
G_{j+1}=G_{j} \cup U_{j} \quad \text { where } \quad U_{j}=G_{j+1} \cap A_{j}
$$

It follows that conditions a), b), c) are satisfied with $c+\varepsilon$ instead of $d$. Finally, the functions

$$
\varphi_{j}=1 /\left(c-\psi_{j+1}\right)+1 /\left(r_{j}^{2}-\left|z-z_{j}\right|^{2}\right), \quad \widetilde{\varphi}_{j}=1 /\left(c-\psi_{j}\right)+1 /\left(r_{j}^{2}-\left|z-z_{j}\right|^{2}\right)
$$

are strongly $q$-convex exhaustions on $U_{j}$ and $G_{j} \cap U_{j}=G_{j} \cap A_{j}$. Let $L$ be an arbitrary compact subset of $G_{j} \cap U_{j}$ and $a=\sup _{L} \psi_{j}<c$. Select $\left.b \in\right] a, c[$ and set

$$
\psi_{j, \eta}=\psi_{j}+\eta \varphi_{j} \quad \text { on } \quad U_{j}, \quad \eta>0
$$

Then $\psi_{j, \eta}$ is an exhaustion of $U_{j}$. As $\varphi_{j}$ is bounded below, we have

$$
L \subset\left\{\psi_{j, \eta}<b\right\} \subset \subset\left\{\psi_{j}<c\right\} \cap U_{j}=G_{j} \cap U_{j}
$$

for $\eta$ small enough. Moreover

$$
(1-\alpha) \psi_{j}+\alpha \psi_{j+1}=\psi-\varepsilon \sum_{0 \leqslant k<j} \theta_{k}-\alpha \varepsilon \theta_{j}
$$

is strongly $q$-convex for all $\alpha \in[0,1]$ and $\varepsilon \leqslant \varepsilon_{0}$ small enough, so Prop. 2.4 implies that $\psi_{j, \eta}$ is strongly $q$-convex. By definition, $G_{j} \cap U_{j}$ is thus $q$-Runge complete in $U_{j}$, and Lemma 4.12 is proved with $X_{c+\varepsilon}$ instead of $X_{d}$. In order to achieve the proof, we consider an increasing sequence $c=c_{0}<c_{1}<\ldots<c_{N}=d$ with $c_{k+1}-c_{k} \leqslant \varepsilon_{0}$ and perform the same construction for each pair $X_{c_{k}} \subset X_{c_{k+1}}$, with $c$ replaced by $c_{k}$ and $\varepsilon=c_{k+1}-c_{k}$.
(4.13) Proposition. For every sublevel set $X_{c} \supset K$, the group $H^{k}\left(X_{c}, \mathscr{Y}\right)$ is Hausdorff and finite dimensional when $k \geqslant q$. Moreover, for $d>c$, the restriction map

$$
H^{k}\left(X_{d}, \mathscr{S}\right) \longrightarrow H^{k}\left(X_{c}, \mathscr{S}\right)
$$

is an isomorphism when $k \geqslant q$ and has a dense range when $k=q-1$.
Proof. Thanks to Lemma 4.12, we are led to consider the restriction maps

$$
\begin{equation*}
H^{k}\left(G_{j+1}, \mathscr{P}\right) \longrightarrow H^{k}\left(G_{j}, \mathscr{P}\right) \tag{4.14}
\end{equation*}
$$

Let us apply the Mayer-Vietoris exact sequence IV-3.11 to $G_{j+1}=G_{j} \cup U_{j}$. For $k \geqslant q$ we have $H^{k}\left(U_{j}, \mathscr{S}\right)=H^{k}\left(G_{j} \cap U_{j}, \mathscr{S}\right)=0$ by Lemma 4.4. Hence we get an exact sequence

$$
\left.\begin{array}{cc}
H^{q-1}\left(G_{j+1}, \mathscr{S}\right) & \longrightarrow H^{q-1}\left(G_{j}, \mathscr{Y}\right) \oplus H^{q-1}\left(U_{j}, \mathscr{Y}\right) \longrightarrow H^{q-1}\left(G_{j} \cap U_{j}, \mathscr{Y}\right) \longrightarrow \\
H^{k}\left(G_{j+1}, \mathscr{P}\right) \longrightarrow & H^{k}\left(G_{j}, \mathscr{Y}\right)
\end{array} \longrightarrow 0 \quad \begin{array}{l}
\longrightarrow
\end{array}\right)
$$

In this sequence, all the arrows are induced by restriction maps, so they define continuous linear operators. We already infer that the map (4.14) is bijective for $k>q$ and surjective for $k=q$. There exist a $\mathscr{S}$-acyclic covering $\mathscr{V}=\left(V_{\alpha}\right)$ of $X_{d}$ and a finite family $\mathscr{V}^{\prime}=$ $\left(V_{\alpha_{1}}^{\prime}, \ldots, V_{\alpha_{p}}^{\prime}\right)$ of open sets such that $V_{\alpha_{j}}^{\prime} \subset \subset V_{\alpha_{j}}$ and $\bigcup V_{\alpha_{j}}^{\prime} \supset \bar{X}_{c}$. Let ${ }^{W} W$ be a locally finite $\mathscr{S}$-acyclic covering of $X_{c}$ which refines $\mathscr{V}^{\prime} \cap X_{c}=\left(V_{\alpha_{j}}^{\prime} \cap X_{c}\right)$. The refinement map

$$
C^{\bullet}(\mathscr{V}, \mathscr{P}) \longrightarrow C^{\bullet}\left(\mathscr{V}^{\prime} \cap X_{c}, \mathscr{S}\right) \longrightarrow C^{\bullet}(\mathscr{W}, \mathscr{S})
$$

is compact because the first arrow is, and it induces a surjective map

$$
H^{k}\left(X_{d}, \mathscr{P}\right) \longrightarrow H^{k}\left(X_{c}, \mathscr{P}\right) \quad \text { for } \quad k \geqslant q
$$

By Schwartz' theorem 1.9, we conclude that $H^{k}\left(X_{c}, \mathscr{P}\right)$ is Hausdorff and finite dimensional for $k \geqslant q$. This is equally true for $H^{q}\left(G_{j}, \mathscr{P}\right)$ because $G_{j}$ is also a global sublevel set $\left\{x \in X ; \psi_{j}(x)<c_{j}\right\}$ containing $K$. Now, the Mayer-Vietoris exact sequence implies that the composite

$$
H^{q-1}\left(U_{j}, \mathscr{S}\right) \longrightarrow H^{q-1}\left(G_{j} \cap U_{j}, \mathscr{S}\right) \xrightarrow{\partial} H^{q}\left(G_{j+1}, \mathscr{Y}\right)
$$

is equal to zero. However, the first arrow has a dense range by Lemma 4.9. As the target space is Hausdorff, the second arrow must be zero; we obtain therefore the injectivity of $H^{q}\left(G_{j+1}, \mathscr{S}\right) \longrightarrow H^{q}\left(G_{j}, \mathscr{S}\right)$ and an exact sequence

$$
\begin{aligned}
H^{q-1}\left(G_{j+1}, \mathscr{S}\right) \longrightarrow H^{q-1}\left(G_{j}, \mathscr{F}\right) \oplus H^{q-1}\left(U_{j}, \mathscr{S}\right) & \longrightarrow H^{q-1}\left(G_{j} \cap U_{j}, \mathscr{S}\right) \longrightarrow 0 \\
g \oplus u & \longmapsto u_{\upharpoonright G_{j} \cap U_{j}}-g_{\left\lceil G_{j} \cap U_{j}\right.} .
\end{aligned}
$$

The argument used in Rem. 1.10 shows that the surjective arrow is open. Let $g \in$ $H^{q-1}\left(G_{j}, \mathscr{S}\right)$ be given. By Lemma 4.9, we can approximate $g_{\left\lceil G_{j} \cap U_{j}\right.}$ by a sequence $u_{\nu \backslash G_{j} \cap U_{j}}, u_{\nu} \in H^{q-1}\left(U_{j}, \mathscr{S}\right)$. Then $w_{\nu}=u_{\nu\left\lceil G_{j} \cap U_{j}\right.}-g_{\left\lceil G_{j} \cap U_{j}\right.}$ tends to zero. As the second map in the exact sequence is open, we can find a sequence

$$
g_{\nu}^{\prime} \oplus u_{\nu}^{\prime} \in H^{q-1}\left(G_{j}, \mathscr{Y}\right) \oplus H^{q-1}\left(U_{j}, \mathscr{Y}\right)
$$

converging to zero which is mapped on $w_{\nu}$. Then $\left(g-g_{\nu}^{\prime}\right) \oplus\left(u_{\nu}-u_{\nu}^{\prime}\right)$ is mapped on zero, and there exists a sequence $f_{\nu} \in H^{q-1}\left(G_{j+1}, \mathscr{S}\right)$ which coincides with $g-g_{\nu}^{\prime}$ on $G_{j}$ and with $u_{\nu}-u_{\nu}^{\prime}$ on $U_{j}$. In particular $f_{\nu\left\lceil G_{j}\right.}$ converges to $g$ and we have shown that

$$
H^{q-1}\left(G_{j+1}, \mathscr{S}\right) \longrightarrow H^{q-1}\left(G_{j}, \mathscr{S}\right)
$$

has a dense range.
Proof of Andreotti-Grauert's Theorem 4.10.. Let 'W be a countable basis of the topology of $X$ consisting of strongly 1-convex open sets $W_{\alpha}$ contained in $\mathscr{G}$-distinguished patches of $X$. Let $L \subset U$ be an arbitrary compact subset. Select a smooth exhaustion function $\psi$ on $X$ such that $\psi$ is strongly $q$-convex on $X \backslash \bar{X}_{b}$ and $L \subset X_{b} \subset \subset U$ for some sublevel set $X_{b}$ of $\psi$; choose $c>b$ such that $X_{c} \subset \subset U$. For every $d \in \mathbb{R}$, we denote by ${ }^{W_{d}} \subset{ }^{`} W$ the collection of sets $W_{\alpha} \in{ }^{W} W$ such that $W_{\alpha} \subset X_{d}$. Then $\mathscr{W}_{d}$ is a $\mathscr{S}$-acyclic covering of $X_{d}$. We consider the sequence of Cech complexes

$$
E_{\nu}^{\bullet}=C^{\bullet}\left(\mathscr{W}_{c+\nu}, \mathscr{P}\right), \quad \nu \in \mathbb{N}
$$

together with the surjective projection maps $E_{\nu+1}^{\bullet} \longrightarrow E_{\nu}^{\bullet}$, and their inverse limit $E^{\bullet}=C^{\bullet}(\mathscr{W}, \mathscr{S})$. Then we have $H^{k}\left(E^{\bullet}\right)=H^{k}(X, \mathscr{Y})$ and $H^{k}\left(E_{\nu}^{\bullet}\right)=H^{k}\left(X_{c+\nu}, \mathscr{Y}\right)$. Propositions 1.11 (a,b,c) and 4.13 imply that $H^{k}(X, \mathscr{Y}) \longrightarrow H^{k}\left(X_{c}, \mathscr{S}\right)$ is bijective for $k \geqslant q$ and has a dense range for $k=q-1$. It already follows that $H^{k}(X, \mathscr{P})$ is Hausdorff for $k \geqslant q$. Now, take an increasing sequence of open sets $X_{c_{\nu}}$ equal to sublevel sets of a sequence of exhaustions $\psi_{\nu}$, such that $U=\bigcup X_{c_{\nu}}$. Then all groups $H^{k}\left(X_{c_{\nu}}, \mathscr{Y}\right)$ are in bijection with $H^{k}(X, \mathscr{S})$ for $k \geqslant q$, and the image of $H^{q-1}\left(X_{c_{\nu+1}}, \mathscr{S}\right)$ in $H^{q-1}\left(X_{c_{\nu}}, \mathscr{S}\right)$ is dense because it contains the image of $H^{q-1}(X, \mathscr{S})$. Proposition 1.11 (a,b,c) again shows that $H^{k}(U, \mathscr{S}) \longrightarrow H^{k}\left(X_{c_{0}}, \mathscr{S}\right)$ is bijective for $k \geqslant q$, and d) shows that $H^{q-1}(X, \mathscr{S}) \longrightarrow H^{q-1}(U, \mathscr{S})$ has a dense range. The theorem follows.

A combination of Andreotti-Grauert's theorem with Th. 3.6 yields the following important consequence.
(4.15) Corollary. Let $\mathscr{S}$ be a coherent sheaf over an analytic scheme ( $X, \mathscr{O}_{X}$ ) with $\operatorname{dim} X \leqslant n$.
a) We have $H^{k}(X, \mathscr{S})=0$ for all $k \geqslant n+1$;
b) If $X$ has no $n$-dimensional compact irreducible component, then we have $H^{n}(X, \mathscr{F})=0$.
c) If $X$ has only finitely many n-dimensional compact irreducible components, then $H^{n}(X, \mathscr{S})$ is finite dimensional.

The special case of 4.15 b ) when $X$ is smooth and $\mathscr{S}$ locally free has been first proved by [Malgrange 1955], and the general case is due to [Siu 1969]. Another consequence is
the following approximation theorem for coherent sheaves over manifolds, which results from Prop. 3.8.
(4.16) Proposition. Let $\mathscr{S}$ be a coherent sheaf over a non compact connected complex manifold $M$ with $\operatorname{dim} M=n$. Let $U \subset M$ be an open subset such that the complement $M \backslash U$ has no compact connected component. Then the restriction map $H^{n-1}(M, \mathscr{P}) \rightarrow$ $H^{n-1}(U, \mathscr{S})$ has a dense range.

## § 5. Grauert's Direct Image Theorem

The goal of this section is to prove the following fundamental result on direct images of coherent analytic sheaves, due to [Grauert 1960].
(5.1) Direct image theorem. Let $X, Y$ be complex analytic schemes and let $F$ : $X \rightarrow Y$ be a proper analytic morphism. If $\mathscr{S}$ is a coherent $\mathbb{O}_{X}$-module, the direct images $R^{q} F_{\star} \mathscr{S}$ are coherent $\mathfrak{O}_{Y}$-modules.

We give below a beautiful proof due to [Kiehl-Verdier 1971], which is much simpler than Grauert's original proof; this proof rests on rather deep properties of nuclear modules over nuclear Fréchet algebras. We first introduce the basic concept of topological tensor product. Our presentation owes much to the seminar lectures by [DouadyVerdier 1973].

## § 5.A. Topological Tensor Products and Nuclear Spaces

The algebra of holomorphic functions on a product space $X \times Y$ is a completion $\mathscr{O}(X) \widehat{\otimes} \mathscr{O}(Y)$ of the algebraic tensor product $\mathscr{O}(X) \otimes \mathscr{O}(Y)$. We are going to describe the construction and the basic properties of the required topological tensor products $\widehat{\otimes}$.

Let $E, F$ be (real or complex) vector spaces equipped with semi-norms $p$ and $q$, respectively. Then $E \otimes F$ can be equipped with any one of the two natural semi-norms $p \otimes_{\pi} q, p \otimes_{\varepsilon} q$ defined by

$$
\begin{aligned}
& p \otimes_{\pi} q(t)=\inf \left\{\sum_{1 \leqslant j \leqslant N} p\left(x_{j}\right) q\left(y_{j}\right) ; t=\sum_{1 \leqslant j \leqslant N} x_{j} \otimes y_{j}, x_{j} \in E, y_{j} \in F\right\}, \\
& p \otimes_{\varepsilon} q(t)=\sup _{\|\xi\|_{p} \leqslant 1,\|\eta\|_{q} \leqslant 1}|\xi \otimes \eta(t)|, \quad \xi \in E^{\prime}, \quad \eta \in F^{\prime} ;
\end{aligned}
$$

the inequalities in the last line mean that $\xi, \eta$ satisfy $|\xi(x)| \leqslant p(x)$ and $|\eta(y)| \leqslant q(y)$ for all $x \in E, y \in F$. Then clearly $p \otimes_{\varepsilon} q \leqslant p \otimes_{\pi} q$, for

$$
p \otimes_{\varepsilon} q\left(\sum x_{j} \otimes y_{j}\right) \leqslant \sum p \otimes_{\varepsilon} q\left(x_{j} \otimes y_{j}\right) \leqslant \sum p\left(x_{j}\right) q\left(y_{j}\right)
$$

Given $x \in E, y \in F$, the Hahn-Banach theorem implies that there exist $\xi, \eta$ such that $\|\xi\|_{p}=\|\eta\|_{q}=1$ with $\xi(x)=p(x)$ and $\eta(y)=q(y)$, hence $p \otimes_{\varepsilon} q(x \otimes y) \geqslant p(x) q(y)$. On the other hand $p \otimes_{\pi} q(x \otimes y) \leqslant p(x) q(y)$, thus

$$
p \otimes_{\varepsilon} q(x \otimes y)=p \otimes_{\pi} q(x \otimes y)=p(x) q(y) .
$$

(5.2) Definition. Let $E, F$ be locally convex topological vector spaces. The topological tensor product $E \widehat{\otimes}_{\pi} F$ (resp. $E \widehat{\otimes}_{\varepsilon} F$ ) is the Hausdorff completion of $E \otimes F$, equipped with the family of semi-norms $p \otimes_{\pi} q$ (resp. $p \otimes_{\varepsilon} q$ ) associated to fundamental families of semi-norms on $E$ and $F$.

Since we may also write

$$
p \otimes_{\pi} q(t)=\inf \left\{\sum\left|\lambda_{j}\right| ; t=\sum \lambda_{j} x_{j} \otimes y_{j}, p\left(x_{j}\right) \leqslant 1, q\left(y_{j}\right) \leqslant 1\right\}
$$

where the $\lambda_{j}$ 's are scalars, we see that the closed unit ball $B\left(p \widehat{\otimes}_{\pi} q\right)$ in $E \widehat{\otimes}_{\pi} F$ is the closed convex hull of $B(p) \otimes B(q)$. From this, we easily infer that the topological dual space $\left(E \widehat{\otimes}_{\pi} F\right)^{\prime}$ is isomorphic to the space of continuous bilinear forms on $E \times F$. Another simple consequence of this interpretation of $B\left(p \widehat{\otimes}_{\pi} q\right)$ is example a) below.

## (5.3) Examples.

a) For all discrete spaces $I, J$, there is an isometry

$$
\ell^{1}(I) \widehat{\otimes}_{\pi} \ell^{1}(J) \simeq \ell^{1}(I \times J) .
$$

b) For Banach spaces $(E, p),(F, q)$, the closed unit ball in $E \widehat{\otimes}_{\varepsilon} F$ is dual to the unit ball $B\left(p^{\prime} \widehat{\otimes}_{\pi} q^{\prime}\right)$ of $E^{\prime} \widehat{\otimes}_{\pi} F^{\prime}$ through the natural pairing extending the algebraic pairing of $E \otimes F$ and $E^{\prime} \otimes F^{\prime}$. If $c_{0}(I)$ denotes the space of bounded sequences on $I$ converging to zero at infinity, we have $c_{0}(I)^{\prime}=\ell^{1}(I)$, hence by duality $c_{0}(I) \widehat{\otimes}_{\varepsilon} c_{0}(J)$ is isometric to $c_{0}(I \times J)$.
c) If $X, Y$ are compact topological spaces and if $C(X), C(Y)$ are their algebras of continuous functions with the sup norm, then

$$
C(X) \widehat{\otimes}_{\varepsilon} C(Y) \simeq C(X \times Y)
$$

Indeed, $C(X)^{\prime}$ is the space of finite Borel measures equipped with the mass norm. Thus for $f \in C(X) \otimes C(Y)$, the $\otimes_{\varepsilon}$-seminorm is given by

$$
\|f\|_{\varepsilon}=\sup _{\|\mu\| \leqslant 1,\|\nu\| \leqslant 1} \mu \otimes \nu(f)=\sup _{X \times Y}|f| ;
$$

the last equality is obtained by taking Dirac measures $\delta_{x}, \delta_{y}$ for $\mu, \nu$ (the inequality $\leqslant$ is obvious). Now $C(X) \otimes C(Y)$ is dense in $C(X \times Y)$ by the Stone-Weierstrass theorem, hence its completion is $C(X \times Y)$, as desired.

Let $f: E_{1} \rightarrow E_{2}$ and $g: F_{1} \rightarrow F_{2}$ be continuous morphisms. For all semi-norms $p_{2}$, $q_{2}$ on $E_{2}, F_{2}$, there exist semi-norms $p_{1}, q_{1}$ on $E_{1}, F_{1}$ and constants $\|f\|=\|f\|_{p_{1}, p_{2}}$, $\|g\|=\|g\|_{q_{1}, q_{2}}$ such that $p_{2} \circ f \leqslant\|f\| p_{1}$ and $q_{2} \circ g \leqslant\|g\| q_{1}$. Then we find

$$
\left(p_{2} \otimes_{\pi} q_{2}\right) \circ(f \otimes g) \leqslant\|f\|\|g\| p_{1} \otimes_{\pi} q_{1}
$$

and a similar formula with $p_{j} \otimes_{\varepsilon} q_{j}$. It follows that there are well defined continuous maps

$$
f \widehat{\otimes}_{\pi} g: E_{1} \widehat{\widehat{\otimes}}_{\pi} F_{1} \longrightarrow E_{2} \widehat{\wedge}_{\pi} F_{2},
$$

$$
f \widehat{\otimes}_{\varepsilon} g: E_{1} \widehat{\otimes}_{\varepsilon} F_{1} \longrightarrow E_{2} \widehat{\otimes}_{\varepsilon} F_{2}
$$

Another simple fact is that $\widehat{\otimes}_{\pi}$ preserves open morphisms:
(5.5) Proposition. If $f: E_{1} \rightarrow E_{2}$ and $g: F_{1} \rightarrow F_{2}$ are epimorphisms, then $f \widehat{\otimes}_{\pi} g$ : $E_{1} \widehat{\otimes}_{\pi} F_{1} \longrightarrow E_{2} \widehat{\otimes}_{\pi} F_{2}$ is an epimorphism.

Proof. Recall that when $E$ is locally convex complete and $F$ Hausdorff, a morphism $u: E \rightarrow F$ is open if and only if $\overline{u(V)}$ is a neighborhood of 0 for every neighborhood of 0 (this can be checked essentially by the same proof as 1.8 b )). Here, for any seminorms $p, q$ on $E_{1}, F_{1}$ the closure of $f \widehat{\otimes}_{\pi} g\left(B\left(p \widehat{\otimes}_{\pi} q\right)\right)$ contains the closed convex hull of $f(B(p)) \otimes g(B(q))$ in which $f(B(p))$ and $g(B(q))$ are neighborhoods of 0 , so it is a neighborhood of 0 in $E \widehat{\otimes}_{\pi} F$.

If $E_{1} \subset E_{2}$ is a closed subspace, every continuous semi-norm $p_{1}$ on $E_{1}$ is the restriction of a continuous semi-norm on $E_{2}$, and every linear form $\xi_{1} \in E_{1}^{\prime}$ such that $\left\|\xi_{1}\right\|_{p_{1}} \leqslant 1$ can be extended to a linear form $\xi_{2} \in E_{2}$ such that $\left\|\xi_{2}\right\|_{p_{2}}=\left\|\xi_{1}\right\|_{p_{1}}$ (Hahn-Banach theorem); similar properties hold for a closed subspace $F_{1} \subset F_{2}$. We infer that

$$
\left(p_{2} \otimes_{\varepsilon} q_{2}\right)_{\left\lceil E_{1} \otimes F_{1}\right.}=p_{1} \otimes_{\varepsilon} q_{1},
$$

thus $E_{1} \widehat{\otimes}_{\varepsilon} F_{1}$ is a closed subspace of $E_{2} \widehat{\otimes}_{\varepsilon} F_{2}$. In other words:
(5.6) Proposition. If $f: E_{1} \rightarrow E_{2}$ and $g: F_{1} \rightarrow F_{2}$ are monomorphisms, then $f \widehat{\otimes}_{\varepsilon} g: E_{1} \widehat{\otimes}_{\varepsilon} F_{1} \longrightarrow E_{2} \widehat{\otimes}_{\varepsilon} F_{2}$ is a monomorphism.

Unfortunately, 5.5 fails for $\widehat{\otimes}_{\varepsilon}$ and 5.6 fails for $\widehat{\otimes}_{\pi}$, even with Fréchet or Banach spaces. It follows that neither $\widehat{\otimes}_{\pi}$ nor $\widehat{\otimes}_{\varepsilon}$ are exact functors in the category of Fréchet spaces. In order to circumvent this difficulty, it is necessary to work in a suitable subcategory.
(5.7) Definition. A morphism $f: E \rightarrow F$ of complete locally convex spaces is said to be nuclear if $f$ can be written as

$$
f(x)=\sum \lambda_{j} \xi_{j}(x) y_{j}
$$

where $\left(\lambda_{j}\right)$ is a sequence of scalars with $\sum\left|\lambda_{j}\right|<+\infty, \xi_{j} \in E^{\prime}$ an equicontinuous sequence of linear forms and $y_{j} \in F$ a bounded sequence.

When $E$ and $F$ are Banach spaces, the space of nuclear morphisms is isomorphic to $E^{\prime} \widehat{\otimes}_{\pi} F$ and the nuclear norm $\|f\|_{\nu}$ is defined to be the norm in this space, namely

$$
\begin{equation*}
\|f\|_{\nu}=\inf \left\{\sum\left|\lambda_{j}\right| ; f=\sum \lambda_{j} \xi_{j} \otimes y_{j},\left\|\xi_{j}\right\| \leqslant 1,\left\|y_{j}\right\| \leqslant 1\right\} . \tag{5.8}
\end{equation*}
$$

For general spaces $E, F$, the equicontinuity of $\left(\xi_{j}\right)$ means that there is a semi-norm $p$ on $E$ and a constant $C$ such that $\left|\xi_{j}(x)\right| \leqslant C p(x)$ for all $j$. Then the definition shows that $f: E \rightarrow F$ is nuclear if and only if $f$ can be factorized as $E \rightarrow E_{1} \rightarrow F_{1} \rightarrow F$ where $E_{1} \rightarrow F_{1}$ is a nuclear morphism of Banach spaces: indeed we need only take $E_{1}$ be equal to the Hausdorff completion $\widehat{E}_{p}$ of $(E, p)$ and let $F_{1}$ be the subspace of $F$ generated by
the closed balanced convex hull of $\left\{y_{j}\right\}$ ( $=$ unit ball in $F_{1}$ ) ; moreover, if $u: S \rightarrow E$ and $v: F \rightarrow T$ are continuous, the nuclearity of $f$ implies the nuclearity of $v \circ f \circ u$; its nuclear decomposition is then $v \circ f \circ u=\sum \lambda_{j}\left(\xi_{j} \circ u\right) \otimes v\left(y_{j}\right)$.
(5.9) Remark. Every nuclear morphism is compact: indeed, we may assume in Def. 5.7 that $\left(y_{j}\right)$ converges to 0 and $\sum\left|\lambda_{j}\right| \leqslant 1$, otherwise we replace $y_{j}$ by $\varepsilon_{j} y_{j}$ with $\varepsilon_{j}$ converging to zero such that $\sum\left|\lambda_{j} / \varepsilon_{j}\right| \leqslant 1$; then, if $U \subset F$ is a neighborhood of 0 such that $\left|\xi_{j}(U)\right| \leqslant 1$ for all $j$, the image $f(U)$ is contained in the closed convex hull of the compact set $\left\{y_{j}\right\} \cup\{0\}$, which is compact.
(5.10) Proposition. If $E, F, G$ are Banach spaces and if $f: E \rightarrow F$ is nuclear, there is a continuous morphism

$$
f \widehat{\otimes} \operatorname{Id}_{G}: E \widehat{\otimes}_{\varepsilon} G \longrightarrow F \widehat{\otimes}_{\pi} G
$$

extending $f \otimes \operatorname{Id}_{G}$, such that $\left\|f \widehat{\otimes} \operatorname{Id}_{G}\right\| \leqslant\|f\|_{\nu}$.
Proof. If $f=\sum \lambda_{j} \xi_{j} \otimes y_{j}$ as in (5.8), then for any $t \in E \otimes G$ we have

$$
\left(f \otimes \operatorname{Id}_{G}\right)(t)=\sum \lambda_{j}\left(\xi_{j} \otimes \operatorname{Id}_{G}(t)\right) \otimes y_{j}
$$

where $\left(\xi_{j} \otimes \operatorname{Id}_{G}\right)(t) \in G$ has norm

$$
\left\|\left(\xi_{j} \otimes \operatorname{Id}_{G}\right)(t)\right\|=\sup _{\eta \in G^{\prime},\|\eta\| \leqslant 1}\left|\eta\left(\xi_{j} \otimes \operatorname{Id}_{G}(t)\right)\right|=\sup _{\eta}\left|\xi_{j} \otimes \eta(t)\right| \leqslant\|t\|_{\varepsilon}
$$

Therefore $\left\|f \otimes \operatorname{Id}_{G}(t)\right\|_{\pi} \leqslant \sum\left|\lambda_{j}\right|\|t\|_{\varepsilon}$, and the infimum over all decompositions of $f$ yields

$$
\left\|f \otimes \operatorname{Id}_{G}(t)\right\|_{\pi} \leqslant\|f\|_{\nu}\|t\|_{\varepsilon}
$$

Proposition 5.10 follows.
If $E$ is a Fréchet space and $\left(p_{j}\right)$ an increasing sequence of semi-norms on $E$ defining the topology of $E$, we have

$$
E=\lim _{\leftrightarrows} \widehat{E}_{p_{j}}
$$

where $\widehat{E}_{p_{j}}$ is the Hausdorff completion of $\left(E, p_{j}\right)$ and $\widehat{E}_{p_{j+1}} \rightarrow \widehat{E}_{p_{j}}$ the canonical morphism. Here $\widehat{E}_{p_{j}}$ is a Banach space for the induced norm $\widehat{p}_{j}$.
(5.11) Definition. A Fréchet space $E$ is said to be nuclear if the topology of $E$ can be defined by an increasing sequence of semi-norms $p_{j}$ such that each canonical morphism

$$
\widehat{E}_{p_{j+1}} \longrightarrow \widehat{E}_{p_{j}}
$$

of Banach spaces is nuclear.

If $E, F$ are arbitrary locally convex spaces, we always have a continuous morphism $E \widehat{\otimes}_{\pi} F \rightarrow E \widehat{\otimes}_{\varepsilon} F$, because $p \otimes_{\varepsilon} q \leqslant p \otimes_{\pi} q$. If $E$, say, is nuclear, this morphism yields in fact an isomorphism $E \widehat{\otimes}_{\varepsilon} F \simeq E \widehat{\otimes}_{\pi} F$ : indeed, by Prop. 5.10, we have $p_{j} \widehat{\otimes}_{\pi} q \leqslant C_{j} p_{j+1} \widehat{\otimes}_{\varepsilon} q$ where $C_{j}$ is the nuclear norm of $\widehat{E}_{p_{j+1}} \rightarrow \widehat{E}_{p_{j}}$. Hence, when $E$
or $F$ is nuclear, we will identify $E \widehat{\otimes}_{\pi} F$ and $E \widehat{\otimes}_{\varepsilon} F$ and omit $\varepsilon$ or $\pi$ in the notation $E \widehat{\otimes} F$.
(5.12) Example. Let $D=\prod D\left(0, R_{j}\right)$ be a polydisk in $\mathbb{C}^{n}$. For any $\left.t \in\right] 0,1[$, we equip $\mathfrak{O}(D)$ with the semi-norm

$$
p_{t}(f)=\sup _{t D}|f| .
$$

The completion of $\left(\mathscr{O}(D), p_{t}\right)$ is the Banach space $E_{t}$ of holomorphic functions on $t D$ which are continuous up to the boundary. We claim that for $t^{\prime}<t<1$ the restriction map

$$
\rho_{t, t^{\prime}}: E_{t^{\prime}} \longrightarrow E_{t}
$$

is nuclear. In fact, for $f \in \mathscr{O}(D)$, we have $f(z)=\sum a_{\alpha} z^{\alpha}$ where $a_{\alpha}=a_{\alpha}(f)$ satisfies the Cauchy inequalities $\left|a_{\alpha}(f)\right| \leqslant p_{t^{\prime}}(f) /\left(t^{\prime} R\right)^{\alpha}$ for all $\alpha \in \mathbb{N}^{n}$. The formula $f=\sum a_{\alpha}(f) e_{\alpha}$ with $e_{\alpha}(z)=z^{\alpha}$ shows that

$$
\left\|\rho_{t, t^{\prime}}\right\|_{\nu} \leqslant \sum\left\|a_{\alpha}\right\|_{p_{t^{\prime}}}\left\|e_{\alpha}\right\|_{p_{t}} \leqslant \sum\left(t^{\prime} R\right)^{-\alpha}(t R)^{\alpha}=\left(1-t / t^{\prime}\right)^{-n}<+\infty .
$$

We infer that $\mathscr{O}(D)$ is a nuclear Fréchet space. It is also in a natural way a fully nuclear Fréchet algebra (see Def. 5.39 below).
(5.13) Proposition. Let $E$ be a nuclear space. A morphism $f: E \rightarrow F$ is nuclear if and only if $f$ admits a factorization $E \rightarrow M \rightarrow F$ through a Banach space $M$.

Proof. By definition, a nuclear map $f: E \rightarrow F$ always has a factorization through a Banach space (even if $E$ is not nuclear). Conversely, if $E$ is nuclear, any continuous linear map $E \rightarrow M$ into a Banach space $M$ is continuous for some semi-norm $p_{j}$ on $E$, so this map has a factorization

$$
E \rightarrow \widehat{E}_{p_{j+1}} \rightarrow \widehat{E}_{p_{j}} \rightarrow M
$$

in which the second arrow is nuclear. Hence any map $E \rightarrow M \rightarrow F$ is nuclear.

## (5.14) Proposition.

a) If $E, F$ are nuclear spaces, then $E \widehat{\otimes} F$ is nuclear.
b) Any closed subspace or quotient space of a nuclear space is nuclear.
c) Any countable product of nuclear spaces is nuclear.
d) Any countable inverse limit of nuclear spaces is nuclear.

Proof. a) If $f: E_{1} \rightarrow F_{1}$ and $g: E_{2} \rightarrow F_{2}$ are nuclear morphisms of Banach spaces, it is easy to check that $f \widehat{\otimes}_{\pi} g$ and $f \widehat{\otimes}_{\varepsilon} g$ are nuclear with $\left\|f \widehat{\otimes}_{?} g\right\|_{\nu} \leqslant\|f\|_{\nu}\|g\|_{\nu}$ in both cases. Property a) follows by applying this to the canonical morphisms $\widehat{E}_{p_{j+1}} \rightarrow \widehat{E}_{p_{j}}$ and $\widehat{F}_{q_{j+1}} \rightarrow \widehat{F}_{q_{j}}$.
c) Let $E_{k}, k \in \mathbb{N}$, be nuclear spaces and $F=\prod E_{k}$. If ( $p_{j}^{k}$ ) is an increasing family of semi-norms on $E_{k}$ as in Def. 5.11, then the topology of $F$ is defined by the family of semi-norms

$$
q_{j}(x)=\max _{0 \leqslant k \leqslant j} p_{j}^{k}\left(x_{k}\right), \quad x=\left(x_{k}\right) \in F .
$$

Then $\widehat{F}_{q_{j}}=\bigoplus_{0 \leqslant k \leqslant j} \widehat{E}_{k, p_{j}^{k}}$ and

$$
\left(\widehat{F}_{q_{j+1}} \rightarrow \widehat{F}_{q_{j}}\right)=\bigoplus_{0 \leqslant k \leqslant j}\left(\widehat{E}_{k, p_{j+1}^{k}} \rightarrow \widehat{E}_{k, p_{j}^{k}}\right) \oplus\left(\widehat{E}_{j+1, p_{j+1}^{j+1}} \rightarrow\{0\}\right)
$$

is easily seen to be nuclear.
b) If $F \subset E$ is closed, then $\widehat{F}_{p_{j}}$ can be identified to a closed subspace of $\widehat{E}_{p_{j}}$, the map $\widehat{F}_{p_{j+1}} \rightarrow \widehat{F}_{p_{j}}$ is the restriction of $\widehat{E}_{p_{j+1}} \rightarrow \widehat{E}_{p_{j}}$ and we have $\widehat{E / F}{ }_{p_{j}} \simeq \widehat{E}_{p_{j}} / \widehat{F}_{p_{j}}$. It is not true in general that the restriction or quotient of a nuclear morphism is nuclear, but this is true for a binuclear $=($ nuclear o nuclear) morphism, as shown by Lemma 5.15 b ) below. Hence $\widehat{F}_{p_{2 j+2}} \rightarrow \widehat{E}_{p_{2 j}}$ and $\widehat{E / F} p_{p_{2 j+2}} \rightarrow \widehat{E / F} F_{p_{2 j}}$ are nuclear, so $\left(p_{2 j}\right)$ is a fundamental family of semi-norms on $F$ or $E / F$, as required in Def. 5.11.
d) follows immediately from b) and c), since $\underset{\longleftarrow}{\lim } E_{k}$ is a closed subspace of $\Pi E_{k}$.
(5.15) Lemma. Let $E, F, G$ be Banach spaces.
a) If $f: E \rightarrow F$ is nuclear, then $f$ can be factorized through a Hilbert space $H$ as a morphism $E \rightarrow H \rightarrow F$.
b) Let $g: F \rightarrow G$ be another nuclear morphism. If $\operatorname{Im}(g \circ f)$ is contained in a closed subspace $T$ of $G$, then $g \circ f: E \rightarrow T$ is nuclear. If $\operatorname{ker}(g \circ f)$ contains a closed subspace $S$ of $E$, the induced map $(g \circ f)^{\sim}: E / S \rightarrow G$ is nuclear.

Proof. a) Write $f=\sum_{j \in I} \xi_{j} \otimes y_{j} \in E^{\prime} \widehat{\otimes}_{\pi} F$ with $\sum\left\|\xi_{j}\right\|\left\|y_{j}\right\|<+\infty$. Without loss of generality, we may suppose $\left\|\xi_{j}\right\|=\left\|y_{j}\right\|$. Then $f$ is the composition

$$
E \longrightarrow \ell^{2}(I) \longrightarrow F, \quad x \longmapsto\left(\xi_{j}(x)\right), \quad\left(\lambda_{j}\right) \longmapsto \sum \lambda_{j} y_{j} .
$$

b) Decompose $g$ into $g=v \circ u$ as in a) and write $g \circ f$ as the composition

$$
E \xrightarrow{f} F \xrightarrow{u} H \xrightarrow{v} G
$$

where $H$ is a Hilbert space. If $\operatorname{Im}(g \circ f) \subset T$ and if $T \subset G$ is closed, then $H_{1}=v^{-1}(T)$ is a closed subspace of $H$ containing $\operatorname{Im}(u \circ f)$. Therefore $g \circ f: E \rightarrow T$ is the composition

$$
E \xrightarrow{f} F \xrightarrow{u} H \xrightarrow{\mathrm{pr}^{\perp}} H_{1} \xrightarrow{v_{\vdash H_{1}}} T
$$

where $f$ is nuclear and $g \circ f: E \longrightarrow T$ is nuclear. Similar proof for $(g \circ f)^{\sim}: E / S \rightarrow G$ by using decompositions $f=v \circ u: E \rightarrow H \rightarrow F$ and

$$
(g \circ f)^{\sim}: E / S \xrightarrow{\widetilde{u}} H / H_{1} \simeq H_{1}^{\perp} \xrightarrow{v_{\mid H_{1}}^{\perp}} F \xrightarrow{g} G
$$

where $H_{1}=\overline{u(S)}$ satisfies $H_{1} \subset \operatorname{ker}(g \circ v) \subset H$.
(5.16) Corollary. Let $E$ be a nuclear space and let $E \rightarrow F$ be a nuclear morphism.
a) If $f(E)$ is contained in a closed subspace $T$ of $F$, then the morphism $f_{1}: E \rightarrow T$ induced by $f$ is nuclear.
b) If $\operatorname{ker} f$ contains a closed subspace $S$ of $E$, then $\tilde{f}: E / S \rightarrow F$ is nuclear.

Proof. Let $E \xrightarrow{u} M \xrightarrow{v} F$ be a factorization of $f$ through a Banach space $M$. In case a), resp. b), $M_{1}=v^{-1}(T)$ is a closed subspace of $M$, resp. $M / \overline{u(S)}$ is a Banach space, and we have factorizations

$$
f_{1}: E \xrightarrow{u_{1}} M_{1} \xrightarrow{v_{1}} T, \quad \tilde{f}: E / S \xrightarrow{\widetilde{u}} M / \overline{u(S)} \xrightarrow{\widetilde{v}} F
$$

where $u_{1}, \widetilde{u}$ are induced by $u$ and $v_{1}, \widetilde{v}$ by $v$. Hence $f_{1}$ and $\tilde{f}$ are nuclear.
(5.17) Proposition. Let $0 \rightarrow E_{1} \rightarrow E_{2} \rightarrow E_{3} \rightarrow 0$ be an exact sequence of Fréchet spaces and let $F$ be a Fréchet space. If $E_{2}$ or $F$ is nuclear, there is an exact sequence

$$
0 \longrightarrow E_{1} \widehat{\otimes} F \longrightarrow E_{2} \widehat{\otimes} F \longrightarrow E_{3} \widehat{\otimes} F \longrightarrow 0
$$

Proof. If $E_{2}$ is nuclear, then so are $E_{1}$ and $E_{3}$ by Prop. 5.14 b). Hence $E_{1} \widehat{\otimes} F \rightarrow E_{2} \widehat{\otimes} F$ is a monomorphism and $E_{2} \widehat{\otimes} F \rightarrow E_{3} \widehat{\otimes} F$ an epimorphism by Prop. 5.6 and 5.5. It only remains to show that

$$
\operatorname{Im}\left(E_{1} \widehat{\otimes} F \longrightarrow E_{2} \widehat{\otimes} F\right)=\operatorname{ker}\left(E_{2} \widehat{\otimes} F \longrightarrow E_{3} \widehat{\otimes} F\right)
$$

and for this, we need only show that the left hand side is dense in the right hand side (we already know it is closed). Let $\varphi \in\left(E_{2} \widehat{\otimes} F\right)^{\prime}$ be a linear form, viewed as a continuous bilinear form on $E_{2} \times F$. If $\varphi$ vanishes on the image of $E_{1} \widehat{\otimes} F$, then $\varphi$ induces a continuous bilinear form on $E_{3} \times F$ by passing to the quotient. Hence $\varphi$ must vanish on the kernel of $E_{2} \widehat{\otimes} F \rightarrow E_{3} \widehat{\otimes} F$, and our density statement follows by the Hahn-Banach theorem.

## § 5.B. Künneth Formula for Coherent Sheaves

As an application of the above general concepts, we now show how topological tensor products can be used to compute holomorphic functions and cohomology of coherent sheaves on product spaces.
(5.18) Proposition. Let $\mathscr{F}$ be a coherent analytic sheaf on a complex analytic scheme $\left(X, \mathscr{O}_{X}\right)$. Then $\mathscr{F}(X)$ is a nuclear space.

Proof. Let $A \subset \Omega \subset \mathbb{C}^{N}$ be an open patch of $X$ such that the image sheaf $\left(i_{A}\right)_{\star} \mathscr{F}_{\wedge A}$ on $\Omega$ has a resolution

$$
\mathscr{O}_{\Omega}^{p_{1}} \longrightarrow \mathscr{O}_{\Omega}^{p_{0}} \longrightarrow\left(i_{A}\right)_{\star} \mathscr{F}_{\mid A} \longrightarrow 0
$$

and let $D \subset \subset \Omega$ be a polydisk. As $D$ is Stein, we get an exact sequence

$$
\begin{equation*}
\mathscr{O}^{p_{1}}(D) \longrightarrow \mathscr{O}^{p_{0}}(D) \longrightarrow \mathscr{F}(A \cap D) \longrightarrow 0 . \tag{5.19}
\end{equation*}
$$

Hence $\mathscr{F}(A \cap D)$ is a quotient of the nuclear space $\mathscr{C}^{p_{0}}(D)$ and so $\mathscr{F}(A \cap D)$ is nuclear by ( 5.14 b ). Let $\left(U_{\alpha}\right)$ be a countable covering of $X$ by open sets of the form $A \cap D$. Then $\mathscr{F}(X)$ is a closed subspace of $\prod \mathscr{F}\left(U_{\alpha}\right)$, thus $\mathscr{F}(X)$ is nuclear by $(5.14 \mathrm{~b}, \mathrm{c})$.
(5.20) Proposition. Let $\mathscr{F}, \mathscr{G}_{\mathcal{G}}$ be coherent sheaves on complex analytic schemes $X, Y$ respectively. Then there is a canonical isomorphism

$$
\mathscr{F} \boxtimes \mathscr{G}(X \times Y) \simeq \mathscr{F}(X) \widehat{\otimes} \mathscr{G}(Y) .
$$

Proof. We show the proposition in several steps of increasing generality.
a) $X=D \subset \mathbb{C}^{n}, Y=D^{\prime} \subset \mathbb{C}^{p}$ are polydisks, $\mathscr{F}=\mathfrak{O}_{X}, \mathscr{G}=\mathfrak{O}_{Y}$.

Let $p_{t}(f)=\sup _{t D}|f|, p_{t}^{\prime}(f)=\sup _{t D^{\prime}}|f|$ and $q_{t}(f)=\sup _{t\left(D \times D^{\prime}\right)}|f|$ be the seminorms defining the topology of $\mathscr{O}(D), \mathscr{O}\left(D^{\prime}\right)$ and $\mathbb{O}\left(D \times D^{\prime}\right)$, respectively. Then $\widehat{E}_{p_{t}}$ is a closed subspace of the space $C(t \bar{D})$ of continuous functions on $t \bar{D}$ with the sup norm, and we have $p_{t} \otimes_{\varepsilon} p_{t}^{\prime}=q_{t}$ by example (5.3 c). Now, $\mathscr{O}(D) \otimes \mathscr{O}\left(D^{\prime}\right)$ is dense in $\mathscr{O}\left(D \times D^{\prime}\right)$, hence its completion with respect to the family $\left(q_{t}\right)$ is $\mathscr{O}(D) \widehat{\otimes}_{\varepsilon} \mathscr{O}\left(D^{\prime}\right)=\mathscr{O}\left(D \times D^{\prime}\right)$.
b) $X$ is embedded in a polydisk $D \subset \mathbb{C}^{n}, X=A \cap D \stackrel{i}{\hookrightarrow} D$,
$i_{\star} \mathscr{F}$ is the cokernel of a morphism $\mathscr{O}_{D}^{p_{1}} \longrightarrow \mathscr{O}_{D}^{p_{0}}$,
$Y=D^{\prime} \subset \mathbb{C}^{p}$ is a polydisk and $\mathscr{G}=\mathscr{O}_{Y}$.
By taking the external tensor product with $\mathfrak{O}_{Y}$, we get an exact sequence

$$
\begin{equation*}
\mathfrak{O}_{D \times Y}^{p_{1}} \longrightarrow \mathscr{O}_{D \times Y}^{p_{0}} \longrightarrow i_{\star} \mathscr{F} \boxtimes \mathscr{O}_{Y} \longrightarrow 0 \tag{5.21}
\end{equation*}
$$

Then we find a commutative diagram

in which the first line is exact as the image of (5.19) by the exact functor $\bullet \widehat{\otimes} \mathscr{O}(Y)$, and the second line is exact because the exact sequence of sheaves (5.21) gives an exact sequence of spaces of sections on the Stein space $D \times Y$; note that $i_{\star} \mathscr{F} \boxtimes ®_{Y}(D \times Y)=$ $\mathscr{F} \boxtimes \mathscr{G}_{Y}(X \times Y)$. As the first two vertical arrows are isomorphisms by a), the third one is also an isomorphism.
c) $X$, $\mathscr{F}$ are as in b),
$Y$ is embedded in a polydisk $D^{\prime} \subset \mathbb{C}^{p}, Y=A^{\prime} \cap D^{\prime} \stackrel{j}{\longrightarrow} D^{\prime}$ and $j_{\star} \mathscr{G}^{G}$ is the cokernel of $\mathfrak{O}_{D^{\prime}}^{q_{1}} \longrightarrow \mathscr{O}_{D^{\prime}}^{q_{0}}$.
Taking the external tensor product with $\mathscr{F}$, we get an exact sequence

$$
\mathscr{F} \boxtimes \mathscr{O}_{D^{\prime}}^{q_{1}} \longrightarrow \mathscr{F} \boxtimes \mathscr{O}_{D^{\prime}}^{q_{0}} \longrightarrow \mathscr{F} \boxtimes j_{\star} \mathscr{G} \longrightarrow 0
$$

and with the same arguments as above we obtain a commutative diagram

d) $X, \mathscr{F}$ are as in b),c) and $Y, \mathscr{C}_{\mathcal{G}}$ are arbitrary.

Then $Y$ can be covered by open sets $U_{\alpha}=A_{\alpha} \cap D_{\alpha}$ embedded in polydisks $D_{\alpha}$, on which the image of $\mathscr{G}$ admits a two-step resolution. We have $\mathscr{F} \boxtimes \mathscr{G}\left(X \times U_{\alpha}\right) \simeq \mathscr{F}(X) \widehat{\otimes}$ $\mathscr{G}_{( }\left(U_{\alpha}\right)$ by c), and the same is true over the intersections $X \times U_{\alpha \beta}$ because $U_{\alpha \beta}=U_{\alpha} \cap U_{\beta}$ can be embedded by the cross product embedding $j_{\alpha} \times j_{\beta}: U_{\alpha \beta} \rightarrow D_{\alpha} \times D_{\beta}$. We have an exact sequence

$$
0 \longrightarrow \mathscr{G}(Y) \longrightarrow \prod_{\alpha} \mathscr{G}\left(U_{\alpha}\right) \longrightarrow \prod_{\alpha, \beta} \mathscr{G}\left(U_{\alpha \beta}\right)
$$

where the last arrow is $\left(c_{\alpha}\right) \mapsto\left(c_{\beta}-c_{\alpha}\right)$, and a commutative diagram with exact lines

$$
\begin{gathered}
0 \longrightarrow \mathscr{F}(X) \widehat{\otimes} \mathscr{Q}(Y) \longrightarrow \prod \mathscr{F}(X) \widehat{\otimes} \mathscr{\otimes} \mathscr{G}\left(U_{\alpha}\right) \longrightarrow \prod \mathscr{F}(X) \widehat{\otimes} \mathscr{G}\left(U_{\alpha \beta}\right) \\
\downarrow \\
0 \longrightarrow \mathscr{F} \boxtimes \mathscr{G}(X \times Y) \longrightarrow \prod \mathscr{F} \boxtimes \mathscr{G}\left(X \times U_{\alpha}\right) \longrightarrow \prod \mathscr{F} \boxtimes \mathscr{G}\left(X \times U_{\alpha \beta}\right) .
\end{gathered}
$$

Therefore the first vertical arrow is an isomorphism.
e) $X, \mathscr{F}, Y, \mathscr{G}$ are arbitrary.

This case is treated exactly in the same way as d) by reversing the roles of $\mathscr{F}, \mathscr{G}$ and by using d) to get the isomorphism in the last two vertical arrows.
(5.22) Corollary. Let $\mathscr{F}, \mathscr{G}$ be coherent sheaves over complex analytic schemes $X, Y$ and let $\pi: X \times Y \rightarrow X$ be the projection. Suppose that $H^{\bullet}\left(Y, \mathcal{G}_{\boldsymbol{G}}\right)$ is Hausdorff.
a) If $X$ is Stein, then $H^{q}(X \times Y, \mathscr{F} \boxtimes \mathscr{G}) \simeq \mathscr{F}(X) \widehat{\otimes} H^{q}(Y, \mathscr{G})$.
b) In general, for every open set $U \subset X$,

$$
\left(R^{q} \pi_{\star}(\mathscr{F} \boxtimes \mathscr{G})\right)(U)=\mathscr{F}(U) \widehat{\otimes} H^{q}(Y, \mathscr{G})
$$

c) If $H^{q}(Y, \mathscr{G})$ is finite dimensional, then

$$
R^{q} \pi_{\star}(\mathscr{F} \boxtimes \mathscr{G})=\mathscr{F} \otimes H^{q}(Y, \mathscr{G}) .
$$

Proof. a) Let $\mathscr{y}=\left(V_{\alpha}\right)$ be a countable Stein covering of $Y$. By the Leray theorem, $H^{\bullet}(Y, \mathscr{G})$ is equal to the cohomology of the Čech complex $C^{\bullet}(\mathscr{V}, \mathscr{G})$. Similarly $X \times \mathscr{V}=$ $\left(X \times V_{\alpha}\right)$ is a Stein covering of $X \times Y$ and we have

$$
H^{q}(X \times Y, \mathscr{F} \boxtimes \mathscr{G})=H^{q}\left(C^{\bullet}(X \times \mathscr{V}, \mathscr{F} \boxtimes \mathscr{G})\right)
$$

However, Prop. 5.20 shows that $C^{\bullet}(X \times \mathcal{V}, \mathscr{F} \boxtimes \mathscr{G})=\mathscr{F}(X) \widehat{\otimes} C^{\bullet}(\mathscr{V}, \mathscr{G})$. Our assumption that $C^{\bullet}(\mathscr{V}, \mathscr{G})$ has Hausdorff cohomology implies that the cocycle and coboundary groups are (nuclear) Fréchet spaces, and that each cohomology group can be computed by means of short exact sequences in this category. By Prop. 5.17, we thus get the desired equality

$$
H^{q}\left(C^{\bullet}(X \times \mathscr{Y}, \mathscr{F} \boxtimes \mathscr{G})\right)=\mathscr{F}(X) \widehat{\otimes} H^{q}\left(C^{\bullet}(\mathscr{V}, \mathscr{G})\right) .
$$

b) The presheaf $U \mapsto \mathscr{F}(U) \widehat{\otimes} H^{q}(Y, \mathscr{G})$ is in fact a sheaf, because the tensor product with the nuclear space $H^{q}(Y, \mathscr{G})$ preserves the exactness of all sequences

$$
0 \longrightarrow \mathscr{F}(U) \longrightarrow \prod \mathscr{F}\left(U_{\alpha}\right) \longrightarrow \prod \mathscr{F}\left(U_{\alpha \beta}\right)
$$

associated to arbitrary coverings $\left(U_{\alpha}\right)$ of $U$. Property b) thus follows from a) and from the fact that $R^{q} \pi_{\star}(\mathscr{F} \boxtimes \mathscr{G})$ is the sheaf associated to the presheaf $U \mapsto H^{q}(U \times Y, \mathscr{F} \boxtimes \mathscr{G})$.
c) is an immediate consequence of b), since the finite dimensionality of $H^{q}(Y, \mathscr{G})$ implies that this space is Hausdorff.
(5.23) Künneth formula. Let $\mathscr{F}, \mathscr{C}_{\mathcal{G}}$ be coherent sheaves over complex analytic schemes $X, Y$ and suppose that the cohomology spaces $H^{\bullet}(X, \mathscr{F})$ and $H^{\bullet}(Y, \mathscr{G})$ are Hausdorff. Then there is an isomorphism

$$
\begin{aligned}
\bigoplus_{p+q=k} H^{p}(X, \mathscr{F}) \widehat{\otimes} H^{q}(Y, \mathscr{G}) & \simeq H^{k}(X \times Y, \mathscr{F} \boxtimes \mathscr{G}) \\
\bigoplus \alpha_{p} \otimes \beta_{q} & \longmapsto \sum \alpha_{p} \smile \beta_{q} .
\end{aligned}
$$

Proof. Consider the Leray spectral sequence associated to the coherent sheaf $\mathscr{S}=\mathscr{F} \boxtimes \mathscr{G}$ and to the projection $\pi: X \times Y \rightarrow X$. By Cor. 5.22 b ) and a use of Čech cohomology, we find

$$
E_{2}^{p, q}=H^{p}\left(X, R^{q} \pi_{\star} \mathscr{F} \boxtimes \mathscr{G}\right)=H^{p}(X, \mathscr{F}) \widehat{\otimes} H^{q}(Y, \mathscr{G}) .
$$

It remains to show that the Leray spectral sequence degenerates in $E_{2}$. For this, we argue as in the proof of Th. IV-15.9. In that proof, we defined a morphism of the double complex $C^{p, q}=\mathscr{F}^{[p]}(X) \otimes \mathscr{G}^{[q]}(Y)$ into the double complex that defines the Leray spectral sequence (in IV-15.9, we only considered the sheaf theoretic external tensor product $\mathscr{F} \boxtimes \mathscr{G}$, but there is an obvious morphism of that one into the analytic tensor product). We get a morphism of spectral sequences which induces at the $E_{2}$-level the obvious morphism

$$
H^{p}(X, \mathscr{F}) \otimes H^{q}(Y, \mathscr{G}) \longrightarrow H^{p}(X, \mathscr{F}) \widehat{\otimes} H^{q}(Y, \mathscr{G}) .
$$

It follows that the Leray spectral sequence $E_{r}^{p, q}$ is obtained for $r \geqslant 2$ by taking the completion of the spectral sequence of $C^{\bullet \bullet \bullet}$. Since this spectral sequence degenerates in $E_{2}$ by the algebraic Künneth theorem, the Leray spectral sequence also satisfies $d_{r}=0$ for $r \geqslant 2$.
(5.24) Remark. If $X$ or $Y$ is compact, the Künneth formula holds with $\otimes$ instead of $\widehat{\otimes}$, and the assumption that both cohomology spaces are Hausdorff is unnecessary. The proof is exactly the same, except that we use ( 5.22 c ) instead of ( 5.22 b ).

## § 5.C. Modules over Nuclear Fréchet Algebras

Throughout this subsection, we work in the category of nuclear Fréchet spaces. Recall that a topological algebra (commutative, with unit element 1) is an algebra $A$ together with a topological vector space structure such that the multiplication $A \times A \rightarrow A$ is
continuous. $A$ is said to be a Fréchet (resp. nuclear) algebra if it is Fréchet (resp. nuclear) as a topological vector space.
(5.25) Definition. A (Fréchet, resp. nuclear) A-module $E$ is a (Fréchet, resp. nuclear) space $E$ with a $A$-module structure such that the multiplication $A \times E \rightarrow E$ is continuous. The module $E$ is said to be nuclearly free if $E$ is of the form $A \widehat{\otimes} V$ where $V$ is a nuclear Fréchet space.

Assume that $A$ is nuclear and let $E$ be a nuclear $A$-module. A nuclearly free resolution $L_{\text {. }}$ of $E$ is an exact sequence of $A$-modules and continuous $A$-linear morphisms

$$
\begin{equation*}
\cdots \longrightarrow L_{q} \xrightarrow{d_{q}} L_{q-1} \longrightarrow \cdots \longrightarrow L_{0} \longrightarrow E \longrightarrow 0 \tag{5.26}
\end{equation*}
$$

in which each $L_{q}$ is a nuclearly free $A$-module. Such a resolution is said to be direct if each $\operatorname{map} d_{q}$ is direct, i.e. if $\operatorname{Im} d_{q}$ has a topological supplementary space in $L_{q-1}$ (as a vector space over $\mathbb{R}$ or $\mathbb{C}$, not necessarily as a $A$-module).
(5.27) Proposition. Every nuclear $A$-module $E$ admits a direct nuclearly free resolution.

Proof. We define the "standard resolution" of $E$ to be

$$
L_{q}=A \widehat{\otimes} \ldots \widehat{\otimes} A \widehat{\otimes} E
$$

where $A$ is repeated $(q+1)$ times; the $A$-module structure of $L_{q}$ is chosen to be the one given by the first factor and we set $d_{0}\left(a_{0} \otimes x\right)=a_{0} x$,

$$
\begin{gathered}
d_{q}\left(a_{0} \otimes \ldots \otimes a_{q} \otimes x\right)=\sum_{0 \leqslant i<q}(-1)^{i} a_{0} \otimes \ldots \otimes a_{i} a_{i+1} \otimes \ldots \otimes a_{q} \otimes x \\
+(-1)^{q} a_{0} \otimes \ldots \otimes a_{q-1} \otimes a_{q} x .
\end{gathered}
$$

Then there is a homotopy operator $h_{q}: L_{q} \rightarrow L_{q+1}$ given by $h_{q}(t)=1 \otimes t$ for all $q\left(h_{q}\right.$, however, is not $A$-linear). This implies easily that $L_{\bullet}$ is a direct nuclearly free resolution.

If $E$ and $F$ are two nuclear $A$-modules, we define $E \widehat{\otimes}_{A} F$ to be

$$
\begin{align*}
& E \widehat{\otimes}_{A} F=\operatorname{coker}(E \widehat{\otimes} A \widehat{\otimes} F \xrightarrow{d} E \widehat{\otimes} F) \quad \text { where }  \tag{5.28}\\
& d(x \otimes a \otimes y)=a x \otimes y-x \otimes a y .
\end{align*}
$$

Then $E \widehat{\otimes}_{A} F$ is a $A$-module which it is not necessarily Hausdorff. If $E \widehat{\otimes}_{A} F$ is Hausdorff, it is in fact a nuclear $A$-module by Prop. 5.14. If $E$ is nuclearly free, say $E=A \widehat{\otimes} V \simeq V \widehat{\otimes} A$, we have $E \widehat{\otimes}_{A} F=V \widehat{\otimes} F$ (which is thus Hausdorff): indeed, there is an exact sequence
$V \widehat{\otimes} A \widehat{\otimes} A \widehat{\otimes} F \longrightarrow V \widehat{\otimes} A \widehat{\otimes} F \longrightarrow V \widehat{\otimes} F \longrightarrow 0$,
$v \otimes a_{0} \otimes a_{1} \otimes x \longmapsto v \otimes a_{0} a_{1} \otimes x-v \otimes a_{0} \otimes a_{1} x, \quad v \otimes a \otimes x \longmapsto v \otimes a x$,
obtained by tensoring the standard resolution of $F$ with $V \widehat{\otimes}$; observe that the tensor product $\widehat{\otimes}$ with a nuclear space preserves exact sequences thanks to Prop. 5.17. We further define $\operatorname{Tô}_{q}^{A}(E, F)$ to be

$$
\begin{equation*}
\operatorname{To}_{q}^{A}(E, F)=H_{q}\left(E \widehat{\otimes}_{A} L_{\bullet}\right) \tag{5.29}
\end{equation*}
$$

where $L_{\bullet}$ is the standard resolution of $F$. There is in fact an isomorphism

$$
\begin{aligned}
E \widehat{\otimes}_{A} L & \stackrel{\simeq}{ } E \widehat{\otimes} A \widehat{\otimes} \cdots \widehat{\otimes} A \widehat{\otimes} F \\
x \otimes_{A}\left(a_{0} \otimes a_{1} \otimes \ldots \otimes a_{q} \otimes y\right) & \longmapsto a_{0} x \otimes a_{1} \otimes \ldots \otimes a_{q} \otimes y
\end{aligned}
$$

where $A$ is repeated $q$ times in the target space. In this isomorphism, the differential becomes

$$
\begin{aligned}
& d_{q}\left(x \otimes a_{1} \otimes \ldots \otimes a_{q} \otimes y\right)=a_{1} x \otimes a_{2} \otimes \ldots \otimes a_{q} \otimes y \\
&+\sum_{1 \leqslant i<q}(-1)^{i} x \otimes a_{1} \otimes \ldots \otimes a_{i} a_{i+1} \otimes \ldots \otimes a_{q} \otimes y \\
&+(-1)^{q} x \otimes a_{1} \otimes \ldots \otimes a_{q-1} \otimes a_{q} y .
\end{aligned}
$$

In particular, we get $\operatorname{Tô}_{0}^{A}(E, F)=E \widehat{\otimes}_{A} F$. Moreover, if we exchange the roles of $E$ and $F$, we obtain a complex which is isomorphic to the above one up to the sign of $d_{q}$, hence $\operatorname{Tô}_{q}^{A}(E, F) \simeq \operatorname{Tô}_{q}^{A}(F, E)$. If $E=A \widehat{\otimes} V$ is nuclearly free, the complex $E \widehat{\otimes}_{A} L_{\bullet}=V \widehat{\otimes} L_{\bullet}$ is exact, thus

$$
E \text { or } F \text { nuclearly free } \Longrightarrow \operatorname{Tô}_{q}^{A}(E, F)=0 \text { for } q \geqslant 1
$$

(5.30) Proposition. For any exact sequence $0 \rightarrow E_{1} \rightarrow E_{2} \rightarrow E_{3} \rightarrow 0$ of nuclear $A$-modules and any nuclear $A$-module $F$, there is an (algebraic) exact sequence

$$
\begin{aligned}
& \cdots \operatorname{Tôr}_{q}^{A}\left(E_{1}, F\right) \longrightarrow \operatorname{Tô}_{q}^{A}\left(E_{2}, F\right) \longrightarrow \operatorname{Tor}_{q}^{A}\left(E_{3}, F\right) \longrightarrow \operatorname{Tôr}_{q-1}^{A}\left(E_{1}, F\right) \cdots \\
& \longrightarrow E_{1} \widehat{\otimes}_{A} F \longrightarrow E_{2} \widehat{\otimes}_{A} F \longrightarrow E_{3} \widehat{\otimes}_{A} F \longrightarrow 0
\end{aligned}
$$

Proof. As the standard resolution $L_{\bullet} \rightarrow F$ is nuclearly free, $L_{q}=A \widehat{\otimes} V_{q}$ say, then $E_{j} \widehat{\otimes}_{A} L_{\bullet}=E_{j} \widehat{\otimes} V_{\bullet}$ for $j=1,2,3$, so we have a short exact sequence of complexes

$$
0 \longrightarrow E_{1} \widehat{\otimes}_{A} L_{\bullet} \longrightarrow E_{2} \widehat{\otimes}_{A} L_{\bullet} \longrightarrow E_{3} \widehat{\otimes}_{A} L_{\bullet} \longrightarrow 0
$$

(5.31) Corollary. For any nuclearly free (possibly non direct) resolution L. of F, there is a canonical isomorphism

$$
\operatorname{Tô}_{q}^{A}(E, F) \simeq H_{q}\left(E \widehat{\otimes}_{A} L_{\bullet}\right)
$$

Proof. Set $B_{q}=\operatorname{Im}\left(L_{q+1} \rightarrow L_{q}\right)$ for all $q \geqslant 0$ and $B_{-1}=F$. Then apply (5.30) to the short exact sequences $0 \rightarrow B_{q} \rightarrow L_{q} \rightarrow B_{q-1} \rightarrow 0$ and the fact that $L_{q}$ is nuclearly free to get

$$
\operatorname{Tô}_{k}^{A}\left(E, B_{q-1}\right) \simeq \begin{cases}\operatorname{Tôr}_{k-1}^{A}\left(E, B_{q}\right) & \text { for } k>1 \\ \operatorname{ker}\left(E \widehat{\otimes}_{A} B_{q} \rightarrow E \widehat{\otimes}_{A} L_{q}\right) & \text { for } k=1\end{cases}
$$

Hence we obtain inductively

$$
\begin{aligned}
\operatorname{Tô}_{q}^{A}(E, F) & =\operatorname{Tôr}_{q}^{A}\left(E, B_{-1}\right) \simeq \ldots \simeq \operatorname{Tô}_{1}^{A}\left(E, B_{q-2}\right) \\
& \simeq \operatorname{ker}\left(E \widehat{\otimes}_{A} B_{q-1} \rightarrow E \widehat{\otimes}_{A} L_{q-1}\right)
\end{aligned}
$$

and a commutative diagram

in which the horizontal line is exact (thanks to the surjectivity of the left oblique arrow and the exactness of the sequence with $E \widehat{\otimes}_{A} B_{q}$ as first term). Therefore $\operatorname{ker}\left(E \widehat{\otimes}_{A}\right.$ $\left.B_{q-1} \rightarrow E \widehat{\otimes}_{A} L_{q-1}\right)$ can be interpreted as the kernel of $E \widehat{\otimes}_{A} L_{q} \rightarrow E \widehat{\otimes}_{A} L_{q-1}$ modulo the image of $E \widehat{\otimes}_{A} L_{q+1} \rightarrow E \widehat{\otimes}_{A} L_{q}$, and this is is precisely the definition of $H_{q}\left(E \widehat{\otimes}_{A} L_{\bullet}\right)$.

Now, we are ready to introduce the crucial concept of transversality.
(5.32) Definition. We say that two nuclear $A$-modules $E, F$ are transverse if $E \widehat{\otimes}_{A} F$ is Hausdorff and if $\operatorname{Tor}_{q}^{A}(E, F)=0$ for $q \geqslant 1$.

For example, a nuclearly free $A$-module $E=A \widehat{\otimes} V$ is transverse to any nuclear $A$-module $F$. Before proving further general properties, we give a fundamental example.
(5.33) Proposition. Let $X, Y$ be Stein spaces and let $U^{\prime} \subset U \subset \subset X, V \subset \subset Y$ be Stein open subsets. If $\mathscr{F}$ is a coherent sheaf over $X \times Y$, then $\mathfrak{O}\left(U^{\prime}\right)$ and $\mathscr{F}(U \times V)$ are transverse over $\mathbb{O}(U)$. Moreover

$$
\mathfrak{O}\left(U^{\prime}\right) \widehat{\otimes}_{\mathscr{O}(U)} \mathscr{F}(U \times V)=\mathscr{F}\left(U^{\prime} \times V\right) .
$$

Proof. Let $\mathscr{L} \bullet \rightarrow \mathscr{F}$ be a free resolution of $\mathscr{F}$ over $U \times V$; such a resolution exists by Cartan's theorem A. Then $\mathscr{L}_{\bullet}(U \times V)$ is a resolution of $\mathscr{F}(U \times V)$ which is nuclearly free over $\mathscr{O}(U)$, for $\mathscr{O}(U \times V)=\mathscr{O}(U) \widehat{\otimes} \mathscr{O}(V)$; in particular, we get

$$
\begin{aligned}
\mathscr{O}\left(U^{\prime}\right) \widehat{\otimes}_{\mathscr{O}(U)} \mathscr{O}(U \times V) & =\mathscr{O}\left(U^{\prime}\right) \widehat{\otimes} \mathscr{O}(V)=\mathscr{O}\left(U^{\prime} \times V\right), \\
\mathscr{O}\left(U^{\prime}\right) \widehat{\otimes}_{\mathscr{O}(U)} \mathscr{L}_{\bullet}(U \times V) & =\mathscr{L}_{\bullet}\left(U^{\prime} \times V\right) .
\end{aligned}
$$

But $\mathscr{L} \bullet\left(U^{\prime} \times V\right)$ is a resolution of $\mathscr{F}\left(U^{\prime} \times V\right)$, so

$$
\operatorname{Tôr}_{q}^{\mathscr{O}(U)}\left(\mathbb{O}\left(U^{\prime}\right), \mathscr{F}(U \times V)\right)= \begin{cases}\mathscr{F}\left(U^{\prime} \times V\right) & \text { for } q=0, \\ 0 & \text { for } q \geqslant 1 .\end{cases}
$$

## (5.34) Properties.

a) If $0 \rightarrow E_{1} \rightarrow E_{2} \rightarrow E_{3} \rightarrow 0$ is an exact sequence of nuclear $A$-modules and if $E_{2}, E_{3}$ are transverse to $F$, then $E_{1}$ is transverse to $F$.
b) Let $A \rightarrow A_{1} \rightarrow A_{2}$ be homomorphisms of nuclear algebras and let $E$ be a nuclear $A$ module. if $A_{1}$ and $A_{2}$ are transverse to $E$ over $A$, then $A_{2}$ is tranverse to $A_{1} \widehat{\otimes}_{A} E$ over $A_{1}$.
c) Let $E^{\bullet}$ be a complex of nuclear A-modules, bounded on the right side, and let $M$ be a nuclear $A$-module which is transverse to all $E^{n}$. If $E^{\bullet}$ is acyclic in degrees $\geqslant k$, then $M \widehat{\otimes}_{A} E^{\bullet}$ is also acyclic in degrees $\geqslant k$.
d) Let $E^{\bullet}, F^{\bullet}$ be complexes of nuclear $A$-modules, bounded on the right side. Let $f^{\bullet}$ : $E^{\bullet} \rightarrow F^{\bullet}$ be a A-linear morphism and let $M$ be a nuclear $A$-module which is transverse to all $E^{q}$ and $F^{q}$. If $f^{\bullet}$ induces an isomorphism $H^{q}\left(f^{\bullet}\right): H^{q}\left(E^{\bullet}\right) \rightarrow H^{q}\left(F^{\bullet}\right)$ in degrees $q \geqslant k$ and an epimorphism in degree $q=k-1$, then

$$
\operatorname{Id}_{M} \widehat{\otimes}_{A} f^{\bullet}: M \widehat{\otimes}_{A} E^{\bullet} \rightarrow M \widehat{\otimes}_{A} F^{\bullet}
$$

has the same property.
Proof. a) is an immediate consequence of the Tôr exact sequence.
To prove b), we need only check that if $A_{1}$ is transverse to $E$ over $A$, then

$$
\operatorname{Tô}_{q}^{A_{1}}\left(A_{2}, A_{1} \widehat{\otimes}_{A} E\right)=\operatorname{Tô}_{q}^{A}\left(A_{2}, E\right), \quad \forall n \geqslant 0
$$

Indeed, if $L_{\bullet}=A \widehat{\otimes} V_{\bullet}$ is a nuclearly free resolution of $E$ over $A$, then $A_{1} \widehat{\otimes}_{A} L_{\mathbf{\bullet}}=$ $A_{1} \widehat{\otimes} V_{\bullet}$ is a nuclearly free resolution of $A_{1} \widehat{\otimes}_{A} E$ over $A_{1}$, since $H_{q}\left(A_{1} \widehat{\otimes}_{A} L_{\mathbf{\bullet}}\right)=$ $\operatorname{Tôr}_{q}^{A}\left(A_{1}, E\right)=0$ for $q \geqslant 1$. Hence

$$
\begin{aligned}
\operatorname{Tôr}_{q}^{A_{1}}\left(A_{2}, A_{1} \widehat{\otimes}_{A} E\right) & =H_{q}\left(A_{2} \widehat{\otimes}_{A_{1}}\left(A_{1} \widehat{\otimes}_{A} L_{\bullet}\right)\right)=H_{q}\left(A_{2} \widehat{\otimes}_{A_{1}}\left(A_{1} \widehat{\otimes} V_{\bullet}\right)\right) \\
& =H_{q}\left(A_{2} \widehat{\otimes} V_{\bullet}\right)=H_{q}\left(A_{2} \widehat{\otimes}_{A} L \bullet\right)=\operatorname{Tô}_{q}^{A}\left(A_{2}, E\right)
\end{aligned}
$$

c) The short exact sequences $0 \rightarrow Z^{q}\left(E^{\bullet}\right) \longleftrightarrow E^{q} \xrightarrow{d^{q}} Z^{q+1}\left(E^{\bullet}\right) \rightarrow 0$ show by backward induction on $q$ that $M$ is transverse to $Z^{q}\left(E^{\bullet}\right)$ for $q \geqslant k-1$. Hence for $q \geqslant k-1$ we obtain an exact sequence

$$
0 \longrightarrow M \widehat{\otimes}_{A} Z^{q}\left(E^{\bullet}\right) \hookrightarrow M \widehat{\otimes}_{A} E^{q} \xrightarrow{d^{q}} M \widehat{\otimes}_{A} Z^{q+1}\left(E^{\bullet}\right) \longrightarrow 0
$$

which gives in particular $Z^{q}\left(M \widehat{\otimes}_{A} E^{\bullet}\right)=B^{q}\left(M \widehat{\otimes}_{A} E^{\bullet}\right)=M \widehat{\otimes}_{A} Z^{q}\left(E^{\bullet}\right)$ for $q \geqslant k$, as desired.
d) is obtained by applying c) to the mapping cylinder $C\left(f^{\bullet}\right)$, as defined in the following lemma (the proof is straightforward and left to the reader).
(5.35) Lemma. If $f^{\bullet}: E^{\bullet} \rightarrow F^{\bullet}$ is a morphism of complexes, the mapping cylinder $C^{\bullet}=C\left(f^{\bullet}\right)$ is the complex defined by $C^{q}=E^{q} \oplus F^{q-1}$ with differential

$$
\left(\begin{array}{cc}
d_{E}^{q} & 0 \\
-f^{q} & d_{F}^{q-1}
\end{array}\right): E^{q} \oplus F^{q-1} \longrightarrow E^{q+1} \oplus F^{q}
$$

Then there is a short exact sequence $0 \rightarrow F^{\bullet-1} \rightarrow C^{\bullet} \rightarrow E^{\bullet} \rightarrow 0$ and the associated connecting homomorphism $\partial^{q}: H^{q}\left(E^{\bullet}\right) \rightarrow H^{q}\left(F^{\bullet}\right)$ is equal to $H^{q}\left(f^{\bullet}\right)$; in particular, $C^{\bullet}$ is acyclic in degree $q$ if and only if $H^{q}\left(f^{\bullet}\right)$ is injective and $H^{q-1}\left(f^{\bullet}\right)$ is surjective.

## §5.D. $A$-Subnuclear Morphisms and Perturbations

We now introduce a notion of nuclearity relatively to an algebra $A$. This notion is needed for example to describe the properties of the $\mathbb{O}(S)$-linear restriction map $\mathbb{O}(S \times$ $U) \rightarrow \mathscr{O}\left(S \times U^{\prime}\right)$ when $U^{\prime} \subset \subset U$.
(5.36) Definition. Let $E$ and $F$ be Fréchet A-modules over a Fréchet algebra $A$ and let $f: E \rightarrow F$ be a A-linear map. We say that
a) $f$ is $A$-nuclear if there exist a scalar sequence $\left(\lambda_{j}\right)$ with $\sum\left|\lambda_{j}\right|<+\infty$, an equicontinuous family of $A$-linears maps $\xi_{j}: E \rightarrow A$ and a bounded sequence $y_{j}$ in $F$ such that for all $x \in E$

$$
f(x)=\sum \lambda_{j} \xi_{j}(x) y_{j} .
$$

b) $f$ is $A$-subnuclear if there exists a Fréchet $A$-module $M$ and an epimorphism $p: M \rightarrow$ $E$ such that $f \circ p$ is $A$-nuclear; if $E$ is nuclear, we also require $M$ to be nuclear.

If $f: E \rightarrow F$ is $A$-nuclear and if $u: S \rightarrow E$ and $v: F \rightarrow T$ are continuous $A$-linear maps then $v \circ f \circ u$ is $A$-nuclear; the same is true for $A$-subnuclear maps. If $V$ and $W$ are nuclear spaces and if $u: V \rightarrow W$ is $\mathbb{C}$-nuclear, then $\operatorname{Id}_{A} \widehat{\otimes} u: A \widehat{\otimes} V \rightarrow A \widehat{\otimes} W$ is $A$-nuclear. From this we infer:
(5.37) Proposition. Let $S, Z$ be Stein spaces and let $U^{\prime} \subset \subset U \subset \subset Z$ be Stein open subsets. Then the restriction $\rho: \mathfrak{O}(S \times U) \rightarrow \mathfrak{O}\left(S \times U^{\prime}\right)$ is $\mathfrak{O}(S)$-nuclear. If $\mathscr{F}$ is a coherent sheaf over $Y \times Z$ with $Y$ Stein and $S \subset \subset Y$, then the restriction map $\rho: \mathscr{F}(S \times U) \rightarrow$ $\mathscr{F}\left(S \times U^{\prime}\right)$ is $\mathbb{O}(S)$-subnuclear.

Proof. As $\mathfrak{O}(S \times U)=\mathscr{O}(S) \widehat{\otimes} \mathscr{O}(U)$ and $\mathfrak{O}(U) \rightarrow \mathscr{O}\left(U^{\prime}\right)$ is $\mathbb{C}$-nuclear, only the second statement needs a proof. By Cartan's theorem A, there exists a free resolution $\mathscr{L} \bullet \rightarrow \mathscr{F}$ over $S \times U$. Then there is a commutative diagram

$$
\begin{array}{cc}
\mathscr{L}_{0}(S \times U) \\
\stackrel{\rho \downarrow}{ } & \mathscr{F}(S \times U) \\
\mathscr{L}_{0}\left(S \times U^{\prime}\right) \longrightarrow \mathscr{F}\left(S \times U^{\prime}\right)
\end{array}
$$

in which the top horizontal arrow is an $\mathscr{O}(S)$-epimorphism and the left vertical arrow is an $\mathscr{O}(S)$-nuclear map; its composition with the bottom horizontal arrow is thus also $\mathfrak{O}(S)$-nuclear.

Let $f: E \rightarrow F$ be a $A$-linear morphism of Fréchet $A$-modules. Suppose that $f(E) \subset$ $F_{1}$ where $F_{1}$ is a closed $A$-submodule of $F$ and let $f_{1}: E \rightarrow F_{1}$ be the map induced by $f$. If $f$ is $A$-nuclear, it is not true in general that $f_{1}$ is $A$-nuclear or $A$-subnuclear, even if $A, E, F$ are nuclear. However:
(5.38) Proposition. With the above notations, suppose $A, E, F$ nuclear. Let $B$ be a nuclear Fréchet algebra and let $\rho: A \rightarrow B$ be a $\mathbb{C}$-nuclear homomorphism. Suppose that $B$ is transverse to $E, F$ and $F / F_{1}$ over $A$. If $f: E \rightarrow F$ is $A$-subnuclear, then $\operatorname{Id}_{B} \widehat{\otimes}_{A} f_{1}: B \widehat{\otimes}_{A} E \rightarrow B \widehat{\otimes}_{A} F_{1}$ is $B$-subnuclear.

Proof. We first show that $\rho \widehat{\otimes}_{A} f_{1}: E=A \widehat{\otimes}_{A} E \rightarrow B \widehat{\otimes}_{A} F_{1}$ is $\mathbb{C}$-nuclear. Since a quotient of a $\mathbb{C}$-nuclear map is $\mathbb{C}$-nuclear by Cor. 5.16 b ), we may suppose for this that $f$ is $A$-nuclear. Write

$$
\begin{array}{rlll}
f(x)=\sum \lambda_{j} \xi_{j}(x) y_{j}, & \xi_{j}: E \rightarrow A, & \sum\left|\lambda_{j}\right|<+\infty, & y_{j} \in F, \\
\rho(t)=\sum \mu_{k} \eta_{k}(t) b_{k}, & \eta_{k}: A \rightarrow \mathbb{C}, & \sum\left|\mu_{k}\right|<+\infty, & b_{k} \in B
\end{array}
$$

as in the definition of $\left(A\right.$-)nuclearity. Then $\rho \widehat{\otimes}_{A} f: E \longrightarrow B \widehat{\otimes}_{A} F$ is $\mathbb{C}$-nuclear: for any $x \in E$, we have $\rho\left(\xi_{j}(x)\right)=\xi_{j}(x) \rho(1)$ in the $A$-module structure of $B$, hence

$$
\begin{aligned}
\rho \widehat{\otimes}_{A} f(x)=\rho \otimes f(1 \otimes x) & =\sum \lambda_{j} \rho\left(\xi_{j}(x)\right) \widehat{\otimes}_{A} y_{j} \\
& =\sum \lambda_{j} \mu_{k}\left(\eta_{k} \circ \xi_{j}\right)(x) b_{k} \widehat{\otimes}_{A} y_{j} .
\end{aligned}
$$

By our transversality assumptions, $B \widehat{\otimes}_{A} F_{1}$ is a closed subspace of $B \widehat{\otimes}_{A} F$. As $\operatorname{Im}\left(\rho \widehat{\otimes}_{A} f\right) \subset B \widehat{\otimes}_{A} F_{1}$, the induced map $\rho \widehat{\otimes}_{A} f_{1}: E \rightarrow B \widehat{\otimes}_{A} F_{1}$ is $\mathbb{C}$-nuclear by Cor. 5.16 a). Finally, there is a commutative diagram

in which the vertical arrows are $B$-linear epimorphisms. The top horizontal arrow is $B$-nuclear by the $\mathbb{C}$-nuclearity of $\rho \widehat{\otimes}_{A} f_{1}$, hence $\operatorname{Id}_{B} \widehat{\otimes}_{A} f_{1}$ is $B$-subnuclear.

Example 5.12 suggests the following definition (which is somewhat less general than some other in current use, but sufficient for our purposes).
(5.39) Definition. We say that a Fréchet algebra $A$ is fully nuclear if the topology of $A$ is defined by an increasing family $\left(p_{t}\right)_{t \in] 0,1[ }$ of multiplicative semi-norms (that is, $p_{t}(x y) \leqslant$ $\left.p_{t}(x) p_{t}(y)\right)$, such that the Banach algebra homomorphism $\widehat{A}_{p_{t^{\prime}}} \rightarrow \widehat{A}_{p_{t}}$ is nuclear for all $t<t^{\prime}<1$.

If $A$ is fully nuclear and $t \in] 0,1]$, we define $A_{t}$ to be the completion of $A$ equipped with the family of semi-norms $\left.p_{\lambda t}, \lambda \in\right] 0,1\left[\right.$. Then $A_{t}$ is again a fully nuclear algebra, and for all $t<t^{\prime}<1$ the canonical map $A_{t^{\prime}} \rightarrow A_{t}$ is nuclear: indeed, for $t \leqslant u<u^{\prime}<t^{\prime}$, there is a factorization

$$
A_{t^{\prime}} \longrightarrow \widehat{A}_{p_{u^{\prime}}} \longrightarrow \widehat{A}_{p_{u}} \longrightarrow A_{t}
$$

If $E$ is a nuclear $A$-module, we say that $E$ is fully $A$-transverse if $E$ is transverse to all $A_{t}$ over $A$. Then by 5.34 b ), each nuclear space

$$
\begin{equation*}
E_{t}=A_{t} \widehat{\otimes}_{A} E \tag{5.40}
\end{equation*}
$$

is a fully $A_{t}$-transverse $A_{t}$-module. If $f: E \rightarrow F$ is a morphism of fully $A$-transverse nuclear modules, there is an induced map

$$
\left.\left.f_{t}=\operatorname{Id}_{A_{t}} \widehat{\otimes}_{A} f: E_{t} \longrightarrow F_{t}, \quad \forall t \in\right] 0,1\right] .
$$

(5.41) Example. Let $X$ be a closed analytic subscheme of an open set $\Omega \subset \mathbb{C}^{N}$, $D=D(a, R) \subset \subset \Omega$ a polydisk and $U=D \cap X$. We have an epimorphism $\mathbb{O}(D) \rightarrow \mathscr{O}(U)$. Denote by $\widetilde{p}_{t}$ the quotient semi-norm of $p_{t}(f)=\sup _{D(a, t R)}|f|$ on $\mathbb{O}(U)$. Then $\mathbb{O}(U)$ equipped with $\left(\widetilde{p}_{t}\right)_{t \in] 0,1[ }$ is a fully nuclear algebra, and $\mathfrak{G}(U)_{t}=\mathscr{O}(D(a, t R) \cap X)$.

Now, let $Y$ be a Stein space, $V \subset \subset Y$ a Stein open subset and $\mathscr{F}$ a coherent sheaf over $X \times Y$. Then Prop. 5.33 shows that $\mathscr{F}(U \times V)$ is a fully transverse nuclear $\mathbb{O}(U)$-module.
(5.42) Subnuclear perturbation theorem. Let $A$ be a fully nuclear algebra, let $E$ and $F$ be two fully $A$-transverse nuclear $A$-modules and let $f, u: E \rightarrow F$ be $A$-linear maps. Suppose that $u$ is $A$-subnuclear and that $f$ is an epimorphism. Then for every $t<1$, the cokernel of

$$
f_{t}-u_{t}: E_{t} \longrightarrow F_{t}
$$

is a finitely generated $A_{t}$-module (as an algebraic module; we do not assert that the cokernel is Hausdorff).

Proof. We argue in several steps. The first step is the following special case.
(5.43) Lemma. Let $B$ be a Banach algebra, $S$ a Fréchet $B$-module and $v: S \rightarrow S$ a $B$-nuclear morphism. Then Coker $\left(\operatorname{Id}_{S}-v\right)$ is a finitely generated $B$-module.

Proof. Let $v(x)=\sum \lambda_{j} \xi_{j}(x) y_{j}$ be a $B$-nuclear decomposition of $v$. We have a factorization

$$
v=\beta \circ \alpha: S \xrightarrow{\alpha} \ell^{1}(B) \xrightarrow{\beta} S
$$

where $\alpha(x)=\left(\lambda_{j} \xi_{j}(x)\right)$ and $\beta\left(t_{j}\right)=\sum t_{j} y_{j}$. Set $w=\alpha \circ \beta: \ell^{1}(B) \rightarrow \ell^{1}(B)$. As $\alpha$ is $B$-nuclear, so is $w$, and $\alpha, \beta$ induce isomorphisms

$$
\operatorname{Coker}\left(\operatorname{Id}_{S}-v\right) \underset{\widetilde{\beta}}{\stackrel{\widetilde{\alpha}}{\rightleftarrows}} \operatorname{Coker}\left(\operatorname{Id}_{\ell^{1}(B)}-w\right)
$$

We are thus reduced to the case when $S$ is a Banach module. Then we write $v=v^{\prime}+v^{\prime \prime}$ with

$$
v^{\prime}(x)=\sum_{1 \leqslant j \leqslant N} \lambda_{j} \xi_{j}(x) y_{j}, \quad v^{\prime \prime}(x)=\sum_{j>N} \lambda_{j} \xi_{j}(x) y_{j} .
$$

For $N$ large enough, we have $\left\|v^{\prime \prime}\right\|<1$, hence $\operatorname{Id}_{S}-v^{\prime \prime}$ is an automorphism and Coker $\left(\operatorname{Id}_{S}-v^{\prime}-v^{\prime \prime}\right)$ is generated by the classes of $y_{1}, \ldots, y_{N}$.

Proof of Theorem 5.42. a) We may suppose that $E$ is nuclearly free and that $u$ is $A$-nuclear, otherwise we replace $f, u$ by their composition with $A \widehat{\otimes} M \longrightarrow M \xrightarrow{p} E$, where $M$ is nuclear and $p: M \rightarrow E$ is an epimorphism such that $u \circ p$ is $A$-nuclear.
b) As in (5.9), there is a $A$-nuclear decomposition $u(x)=\sum \lambda_{j} \xi_{j}(x) y_{j}$ where $\left(y_{j}\right)$ converges to 0 in $F$. Since $f$ is an epimorphism, we can find a sequence $\left(x_{j}\right)$ converging to 0 in $E$ such that $f\left(x_{j}\right)=y_{j}$. Hence we have $u=f \circ v$ where $v(x)=\sum \lambda_{j} \xi_{j}(x) x_{j}$ is a $A$-nuclear endomorphism of $E$, and the cokernel of $f-u$ is the image by $f$ of the cokernel of $\operatorname{Id}_{E}-v$.
c) By a), b) we may suppose that $F=E=A \widehat{\otimes} M, f=\operatorname{Id}_{E}$ and that $u$ is $A$-nuclear. Let $B$ be the Banach algebra $B=\widehat{A}_{p_{t}}$. Then $B \widehat{\otimes}_{A} E=B \widehat{\otimes} M$ is a Fréchet $B$-module and
$\operatorname{Id}_{B} \widehat{\otimes}_{A} u$ is $B$-nuclear. By Lemma 5.42, $\operatorname{Id}_{B} \widehat{\otimes}_{A} \operatorname{Id}_{E}-\operatorname{Id}_{B} \widehat{\otimes}_{A} u$ has a finitely generated cokernel over $B$. Now, there is an obvious morphism $B \rightarrow A_{t}$, hence by taking the tensor product with $A_{t} \widehat{\otimes}_{B} \bullet$ we get

$$
A_{t} \widehat{\otimes}_{B}\left(B \widehat{\otimes}_{A} E\right)=A_{t} \widehat{\otimes}_{B}(B \widehat{\otimes} M)=A_{t} \widehat{\otimes} M=A_{t} \widehat{\otimes}_{A} E=E_{t}
$$

and we see that

$$
\operatorname{Id}_{E_{t}}-u_{t}=\operatorname{Id}_{A_{t}} \widehat{\otimes}_{A} \operatorname{Id}_{E}-\operatorname{Id}_{A_{t}} \widehat{\otimes}_{A} u
$$

has a finitely generated cokernel over $A_{t}$.

## § 5.E. Proof of the Direct Image Theorem

We first prove a functional analytic version of the result, which appears as a relative version of the Cartan-Serre/Schwartz theorem (Th. 1.9).
(5.44) Theorem. Let $A$ be a fully nuclear algebra, $E^{\bullet}$ and $F^{\bullet}$ complexes of fully $A$ transverse nuclear $A$-modules. Let $f^{\bullet}: E^{\bullet} \rightarrow F^{\bullet}$ be a morphism of complexes such that each $f^{q}$ is $A$-subnuclear. Suppose that $E^{\bullet}$ and $F^{\bullet}$ are bounded on the right and that $H^{q}\left(f^{\bullet}\right)$ is an isomorphism for each $q$. Then for every $t<1$, there is a complex $L^{\bullet}$ of finitely generated free $A_{t}$-modules and a complex morphism $h^{\bullet}: L^{\bullet} \rightarrow E_{t}^{\bullet}$ which induces an isomorphism on cohomology.

Proof. a) We first show the following statement:
Suppose that $E_{t}^{\bullet}$ and $F_{t}^{\bullet}$ are acyclic in degrees $>q$. Then for every $t^{\prime}<t$, the cohomology space $H^{q}\left(E_{t^{\prime}}^{\bullet}\right) \simeq H^{q}\left(F_{t^{\prime}}^{\bullet}\right)$ is a finitely generated $A_{t^{\prime}}$-module.
Indeed, the exact sequences $0 \rightarrow Z^{k}\left(E_{t}^{\bullet}\right) \rightarrow E_{t}^{k} \rightarrow Z^{k+1}\left(E_{t}^{\bullet}\right) \rightarrow 0$ show by backward induction on $k$ that $Z^{k}\left(E_{t}^{\bullet}\right)$ is fully $A_{t}$-transverse for $k \geqslant q$. The same is true for $Z^{k}\left(F_{t}^{\bullet}\right)$. Then $f_{t}^{q}$ is a $A_{t}$-subnuclear map from $Z^{q}\left(E_{t}^{\bullet}\right)$ into $F_{t}^{q}$, and its image is contained in the closed subspace $Z^{q}\left(F_{t}^{\bullet}\right)$. By Prop. 5.38 , for all $t^{\prime \prime}<t$, the map $f_{t^{\prime \prime}}^{q}=\operatorname{Id}_{A_{t^{\prime \prime}}} \widehat{\otimes}_{A_{t}} f_{t}^{q}$ is a $A_{t^{\prime \prime}-\text { subnuclear map }} Z^{q}\left(E_{t^{\prime \prime}}^{\bullet}\right) \rightarrow Z^{q}\left(F_{t^{\prime \prime}}^{\bullet}\right)$. By Prop. 5.34 d$), H^{\bullet}\left(f_{t^{\prime \prime}}^{\bullet}\right)$ is an isomorphism in all degrees, hence

$$
d_{t^{\prime \prime}}^{q} \oplus f_{t^{\prime \prime}}^{q}: F_{t^{\prime \prime}}^{q-1} \oplus Z^{q}\left(E_{t^{\prime \prime}}^{\bullet}\right) \longrightarrow Z^{q}\left(F_{t^{\prime \prime}}^{\bullet}\right)
$$

is surjective. By the subnuclear perturbation theorem, the map

$$
d_{t^{\prime}}^{q} \oplus 0=\operatorname{Id}_{A_{t^{\prime}}} \widehat{\otimes}_{A_{t^{\prime \prime}}}\left(\left(d_{t^{\prime \prime}}^{q} \oplus f_{t^{\prime \prime}}^{q}\right)-\left(0 \oplus f_{t^{\prime \prime}}^{q}\right)\right)
$$

has a finitely generated $A_{t^{\prime}}$-cokernel for $t^{\prime}<t^{\prime \prime}<t$, as desired.
b) Let $N$ be an index such that $E^{k}=F^{k}=0$ for $k>N$. Fix a sequence $t<\ldots<t_{q}<$ $t_{q+1}<\ldots<t_{N}<1$. To prove the theorem, we construct by backward induction on $q$ a finitely generated free module $L^{q}$ over $A_{t_{q}}$ and morphisms $d^{q}: L^{q} \rightarrow L_{t_{q}}^{q+1}, h^{q}: L^{q} \rightarrow E_{t_{q}}^{q}$ such that
i) $L_{\geqslant q, t_{q}}^{\bullet}: 0 \rightarrow L^{q} \rightarrow L_{t_{q}}^{q+1} \rightarrow \cdots \rightarrow L_{t_{q}}^{N} \rightarrow 0$ is a complex and $h_{\geqslant q, t_{q}}^{\bullet}: L_{\geqslant q, t_{q}}^{\bullet} \rightarrow E_{t_{q}}^{\bullet}$ is a complex morphism.
ii) The mapping cylinder $M_{q}^{\bullet}=C\left(h_{\geqslant q, t_{q}}^{\bullet}\right)$ defined by $M_{q}^{k}=\bigoplus_{k \in \mathbb{Z}}\left(L_{\geqslant q, t_{q}}^{k} \oplus E_{t_{q}}^{k-1}\right)$ is acyclic in degrees $k>q$.

Suppose that $L^{k}$ has been constructed for $k \geqslant q$. Consider the mapping cylinder $N_{q}^{\bullet}=$ $C\left(f_{\dot{t}_{q}}^{\bullet} \circ h_{\geqslant q, t_{q}}^{\bullet}\right)$ and the complex morphism

$$
M_{q}^{\bullet} \longrightarrow N_{q}^{\bullet}, \quad L_{\geqslant q, t_{q}}^{k} \oplus E_{t_{q}}^{k-1} \longrightarrow L_{\geqslant q, t_{q}}^{k} \oplus F_{t_{q}}^{k-1}
$$

given by $\operatorname{Id} \oplus f_{t_{q}}^{k-1}$. This morphism is $A_{t_{q}}$-subnuclear in each degree and induces an isomorphism in cohomology (compare the cohomology of the short exact sequences associated to each mapping cylinder, with the obvious morphism between them). Moreover, $M_{q}^{\bullet}$ and $N_{q}^{\bullet}$ are acyclic in degrees $k>q$. By step a), the cohomology space $H^{q}\left(M_{q, t_{q-1}}^{\bullet}\right)$ is a finitely generated $A_{t_{q-1}}$-module. Therefore, we can find a finitely generated free $A_{t_{q-1}}$-module $L^{q-1}$ and a morphism

$$
d^{q-1} \oplus h^{q-1}: L^{q-1} \rightarrow M_{q, t_{q-1}}^{q}=L_{t_{q-1}}^{q} \oplus E_{t_{q-1}}^{q-1}
$$

such that the image is contained in $Z^{q}\left(M_{q, t_{q-1}}^{\bullet}\right)$ and generates the cohomology space $H^{q}\left(M_{q, t_{q-1}}^{\bullet}\right)$. As $M_{q, t_{q-1}}^{q-1}=E_{t_{q-1}}^{q-2}$, this means that $M_{q-1}^{\bullet}$ is also acyclic in degree $q$. Thus $L^{q-1}$, together with the maps ( $d^{q-1}, h^{q-1}$ ) satisfies the induction hypotheses for $q-1$, and $L_{t}^{\bullet}$ together with the induced map $h_{t}^{\bullet}: L_{t}^{\bullet} \rightarrow E_{t}^{\bullet \bullet}$ is the required morphism of complexes.

Proof of Theorem 5.1. Let $X, Y$ be complex analytic schemes, let $F: X \rightarrow Y$ be a proper analytic morphism and let $\mathscr{S}$ be a coherent sheaf over $X$. Fix a point $y_{0} \in Y$, a neighborhood of $y_{0}$ which is isomorphic to a closed analytic subscheme of a Stein open set $W \subset \mathbb{C}^{n}$ and a polydisk $D^{0}=D\left(y_{0}, R_{0}\right) \subset \subset W$. The compact set $K=F^{-1}\left(\bar{D}^{0} \cap Y\right)$ can be covered by finitely many open subsets $U_{\alpha}^{0} \subset \subset X$ which possess embeddings as closed analytic subschemes of Stein open sets $\Omega_{\alpha}^{0} \subset \mathbb{C}^{N_{\alpha}}$. Let $\Omega_{\alpha}^{\prime} \subset \subset \Omega_{\alpha} \subset \subset \Omega_{\alpha}^{0}$ be Stein open subsets such that $U_{\alpha}=U_{\alpha}^{0} \cap \Omega_{\alpha}$ and $U_{\alpha}^{\prime}=U_{\alpha}^{0} \cap \Omega_{\alpha}^{\prime}$ still cover $K$. Let $i_{\alpha}: U_{\alpha}^{0} \rightarrow \Omega_{\alpha}^{0}$ and $j: Y \cap D^{0} \rightarrow D^{0}$ be the embeddings and $\mathscr{S}_{\alpha}=\left(i_{\alpha} \times(j \circ F)\right)_{\star} \mathscr{S}$ the image sheaf of $\mathscr{S}$ on $\Omega_{\alpha}^{0} \times D^{0}$. Let $D \subset \subset D^{0}$ be a concentric polydisk. Then $\mathscr{S}\left(U_{\alpha} \cap F^{-1}(D)\right)=\mathscr{S}_{\alpha}\left(\Omega_{\alpha} \times D\right)$ is a fully transverse $\mathscr{O}(D)$-module by Ex. 5.41, and so is $\mathscr{S}\left(U_{\alpha}^{\prime} \cap F^{-1}(D)\right)=\mathscr{S}_{\alpha}\left(\Omega_{\alpha}^{\prime} \times D\right)$. Moreover, the restriction map

$$
\mathscr{P}\left(U_{\alpha} \cap F^{-1}(D)\right) \longrightarrow \mathscr{S}\left(U_{\alpha}^{\prime} \cap F^{-1}(D)\right)
$$

is $\mathfrak{G}(D)$-subnuclear by Prop. 5.37 applied to $\mathscr{F}=\mathscr{S}_{\alpha}$. For every Stein open set $V \subset D$, Prop. 5.33 shows that

$$
\mathscr{O}(V) \widehat{\otimes}_{\mathscr{O}(D)} \mathscr{P}\left(U_{\alpha} \cap F^{-1}(D)\right)=\mathscr{S}\left(U_{\alpha} \cap F^{-1}(V)\right)
$$

Denote by $\mathscr{U} \cap F^{-1}(D)$ the collection $\left(U_{\alpha} \cap F^{-1}(D)\right)$ and use a similar notation with $U^{\prime}=\left(U_{\alpha}^{\prime}\right)$. As $\mathscr{U} \cap F^{-1}(D), \mathscr{U}^{\prime} \cap F^{-1}(D)$ are Stein coverings of $F^{-1}(D)$, the Leray theorem applied to the alternate Čech complex of $\mathscr{S}$ over $\mathscr{U} \cap F^{-1}(D)$ and $\mathscr{U}^{\prime} \cap F^{-1}(D)$ gives an isomorphism

$$
H^{\bullet}\left(A C^{\bullet}\left(\mathscr{U} \cap F^{-1}(D), \mathscr{S}\right)\right)=H^{\bullet}\left(A C^{\bullet}\left(U^{\prime} \cap F^{-1}(D), \mathscr{S}\right)\right)=H^{\bullet}\left(F^{-1}(D), \mathscr{S}\right) .
$$

By the above discussion, $A C^{\bullet}\left(\mathscr{U} \cap F^{-1}(D), \mathscr{S}\right)$ and $A C^{\bullet}\left(U^{\prime} \cap F^{-1}(D), \mathscr{S}\right)$ are finite complexes of fully transverse nuclear $\mathfrak{O}(D)$-modules, the restriction map

$$
A C^{\bullet}\left(\mathscr{U} \cap F^{-1}(D), \mathscr{S}\right) \longrightarrow A C^{\bullet}\left(U^{\prime} \cap F^{-1}(D), \mathscr{S}\right)
$$

is $\mathfrak{O}(D)$-subnuclear and induces an isomorphism on cohomology groups. Set $D=D\left(y_{0}, R\right)$ and $D_{t}=D\left(y_{0}, t R\right)$. Theorem 5.44 shows that for every $t<1$ there is a complex of finitely generated free $\mathfrak{O}$-modules $\mathscr{L} \bullet$ and a $\mathbb{O}\left(D_{t}\right)$-linear morphism of complexes

$$
\mathscr{L} \bullet\left(D_{t}\right) \rightarrow A C^{\bullet}\left(\mathscr{U} \cap F^{-1}\left(D_{t}\right), \mathscr{S}\right)
$$

which induces an isomorphism on cohomology. Let $V \subset D_{t}$ be an arbitrary Stein open set. By Prop. 5.34 d$)$ applied with $M=\mathscr{O}(V)$, we conclude that $\mathscr{L} \bullet(V) \rightarrow A C^{\bullet}(U \cap$ $\left.F^{-1}(V), \mathscr{S}\right)$ induces an isomorphism on cohomology. If we take the direct limit as $V$ runs over all Stein neighborhoods of a point $y \in Y \cap D_{t}$, we see that $\mathscr{H}^{q}\left(\mathscr{L}^{\bullet}\right) \simeq R^{q} F_{\star} \mathscr{S}$ over $Y \cap D_{t}$, hence $R^{q} F_{\star} \mathscr{S}$ is $\mathscr{O}_{Y}$-coherent near $y_{0}$.

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