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## PREFACE

This book is a revised and expanded version of a series of talks given in Hanoi at the Viện Toán học (Mathematical Institute) in July, 1978. The purpose of the book is the same as the purpose of the talks: to make certain recent applications of p-adic analysis to number theory accessible to graduate students and researchers in related fields. The emphasis is on new results and conjectures, or new interpretations of earlier results, which have come to light in the past couple of years and which indicate intriguing and as yet mperfectly understood new connections between algebraic number theory, algebraic geometry, and p-adic analysis.
occasionally state without proof or assume some familiarity with facts or techniques of other fields: algebraic geometry (Chapter III), algebraic number theory (Chapter IV), analysis (the Appendix). But I include down-to-earth examples and words of otivation whenever possible, so that even a reader with little background in these areas should be able to see what's going on

Chapter I.contains the basic information about p-adic numbers and p-adic analysis needed for what follows. Chapter II describes the construction and properties of p-adic Dirichlet L-functions, ncluding Leopoldt's formula for the value at 1 , using the approach of p-adic integration. The p-adic gamma function and log gamma function are introduced, their properties are developed and compared with the identities satisfied by the classical gamma function, and two formulas relating them to the p-adic L-functions $L_{p}(s, X)$ are proved. The first formula--expressing $L_{p}^{\prime}(0, \chi)$ in terms of special
values of log gamma--will be used later (Chapter IV) in the discussion of Gross' p-adic regulator. The other formula--a p-adic Stirling series for log gamma near infinity--will be a key motivating example for the p-adic Stieltjes transform, discussed in the Appendix.

Chapter III is devoted primarily to proving a p-adic formula for Gauss sums, which expresses them essentially as values of the p-adic gama function. The approach emphasizes the analogy with the complex-analytic periods of differentials on certain special curves, and uses some algebraic geometry. The reader who is interested in a treatment that is more "elementary" and self-contained (but more computational rather than geometric) is referred to [62].

Chapter IV discusses two different types of p-adic regulators. One, due to Leopoldt, is connected with the behavior of $L_{p}(s, X)$ at $s=1$; the other, due to Gross, is connected with the behavior at $s=0$. Conjectures describing these connections between regulators and L -functions are explained and compared to the classical case. The conjectures are proved in the case of a one-dimensional character $X$ with base field $Q$ (the "abelian over $Q$ " case). The proof of Gross' conjecture in this case combines the formula for $L_{p}^{\mathrm{p}}(0, \mathrm{X})$ in Chapter II and the p -adic formula for Gauss sums in Chapter III, together with a p-adic version of the linear independence over $\bar{Q}$ of logarithms of algebraic numbers (Baker's theoren). This proof provides the culmination of the main part of the book.

The Appendix concerns some general constructions in p-adic analysis: the Stieltjes transform and the Shnirelman integral. I first use the Stieltjes transform to highlight the analogy between the p-adic and classical log gama functions. I then give a complete account of M. M. Vishik's p-adic spectral theorem. This material has been relegated to the Appendix because it has not yet led to new number theoretic or algebra-geometric facts, perhaps because Vishik's theory is not very well known.

I would like to thank N. M. Katz, whose Spring 1978 lectures at

Princeton provided the explanations of the algebraic geometry and p-adic cohomology given in Chapter III; R. Greenberg, whose seminar talks at the University of Washington in October 1979 and whose comments on the manuscript were of great help in writing Chapter IV; B. H. Gross, whose preprint [35] and correspondence were the basis for the second half of Chapter IV; and M. M. Vishik, whose preprint [95] is given in modified form in $5 \$ 3-4$ of the Appendix.

I am also grateful to Ju. I. Manin and A. A. Kirillov for the stimulation provided by their seminars on Diophantine geometry and p-adic analysis during my stays in Moscow in 1974-75 and in Spring 1978; and to the Vietnamese mathematicians, in particular Lê-vănThiêm, Hà-huy-Khoâi, Vtrơng-ngọc-Châu and Dỗ-ngọc-Dị̣̂p, for their hospitality, which contributed to a fruitful and enjoyable visit to Hanoi.

Seattle
Neal Koblitz
April 1980

FRONTISPIECE: Artist's conception of the construction of the 2-adic number system as an inverse limit. By Professor A. T. Fomenko of Moscow State University.
I. BASICS

In some places in this chapter detailed proofs and computations are omitted, in order not to bore the reader before we get to the main subject matter. These details are readily available (see, for example, [53]).

1. History (very brief?

| Kummer and Hense1 | $\begin{array}{r} 1850- \\ 1900 \end{array}$ | introduced p-adic numbers and developed their basic properties |
| :---: | :---: | :---: |
| Minkowski. | 1884 | proved: an equation $a_{1} x_{1}^{2}+\ldots+a_{n} x_{n}^{2}=0$ ( $a_{i}$ rational) is solvable in the rational numbers if and only if it is solvable in the reals and in the p-adic numbers for all primes $p$ (see [13,84]) |
| Tate. | 1950 | Fourier analysis on p-adic groups; pointed toward interrelations between p-adic numbers and L-functions and representation theory (see [59]) |
| Dwork.. | 1960 | used p-adic analysis to prove rationality of the zeta-function of an algebraic variety defined over a finite field, part of the Weil conjectures (see $[25,53]$ ) |
| (Kummer....... | 1851 | congruences for Bernoulli numbers--but he approached them in an ad hoc way, without p-adic numbers) |
| Kubota-Leopoldt | 1964 | interpretation of Kummer congruences for Bernoulli numbers using p-adic zeta-function |
| Iwasawa, Serre, Mazur, Manin, Katz, others | past <br> 15 <br> years | p -adic theories for many arithmetically interesting functions |

Dwork, past p-adic differential equations, p-adic their students years

## 2. Basic concepts

Let $p$ be a prime number, fixed once and for all. The "p-adic numbers" are all expressions of the form

$$
\begin{aligned}
& a_{m} p^{m}+a_{m+1} p^{m+1}+a_{m+2} p^{m+2}+\ldots, \\
& \text { re the } a_{i} \in\{0,1,2, \ldots, p-1\} \text { are digits, and } m \text { is any }
\end{aligned}
$$

integer. These expressions form a field $(+$ and $\times$ are defined in the obvious way), which contains the nonnegative integers

$$
n=a_{0}+a_{1} p+\ldots+a_{r} p^{r} \text { ("n written to the base } p^{\prime \prime} \text { ), }
$$

and hence contains the field of rational numbers $Q$. For example,

$$
-1=(p-1)+(p-1)_{p}+(p-1) p^{2}+\ldots
$$

$$
\frac{-a_{0}}{p-1}=a_{0}+a_{0} p+a_{0} p^{2}+\ldots
$$

as is readily seen by adding 1 to the first expression on the right and multiplying the second expression on the right by 1-p.

An equivalent way to define the field $Q_{p}$ of $p$-adic numbers is as the completion of $Q$ under the " $p$-adic metric" determined by the norm | | $\mathrm{p}: \mathrm{Q} \rightarrow$ nonnegative real numbers, defined by

$$
\left|\frac{2}{b}\right|_{p}=p^{\text {ord }_{p} b^{b}-\text { ord }_{p}^{a}}, \quad|0|_{p}=0
$$

where ord ${ }_{p}$ of a nonzero integer is the highest power of $p$ dividing it. Under this norm, numbers highly divisible by $p$ are "smail", while numbers with $p$ in the denominator are "large". For example, $|250|_{5}=1 / 125, \quad|1 / 250|_{5}=125$. Clearly, $\left|\left.\right|_{p}\right.$ is multiplicative, because ord $p$ behaves like log:

$$
\operatorname{ord}_{p}(x y)=\operatorname{ord}_{p} x+\operatorname{ord}_{p} y
$$

Also note that $|\mathrm{n}|_{\mathrm{p}} \leq 1$ for n an integer.
It is not hard to verify that the completion of $Q$ under the $p$-adic metric can be identified with the set $Q_{p}$ of " $p$-adic expansions" $a_{m} p^{m}+a_{m+1} p^{m+1}+\ldots$. The norm $\left.\right|_{p}$ is easy to
evaluate on an element of $Q_{p}$ written in its p adic expansion: if $x=a_{\mathrm{m}^{2}} \mathrm{p}^{m}+a_{m+1} p^{m+1}+\ldots$ with $a_{m} \neq 0$, then $|x|_{p}=p^{-m}$.

Thus, $Q_{p}$ is obtained from $\left|\left.\right|_{p}\right.$ in the same way as the real number field R is obtained from the usual absolute value $\mid$ |: as the completion of $Q$. In fact, a theorem of Ostrowski (see [13] or [53]) says that any norm on $Q$ is equivalent to the usual | | or to $\left|\left.\right|_{p}\right.$ for some $p$. Hence, together with $R$, the various $Q_{p}$ make up all possible completions of $Q$ :

$$
\begin{array}{llllll}
\mathrm{R} & Q_{2} & Q_{3} & Q_{5} & \cdots & Q_{p} \\
& \supset & & & \\
& & &
\end{array}
$$

Q

Oftern, a situation can be studied more easily over $R$ and $Q_{p}$ than over $Q$; and then the information obtained can be put together to conclude something about the situation over Q. For example, one can readily show that a rational number has a square root in $Q$ if and only if it has a square root in $R$ and for all $p$ has a square root in $Q_{p}$. This assertion is a special case of the Hasse-Minkowski theorem (see $\S 1$ above).

In addition to multiplicativity, the other basic property of a norm $|\mid$ on a field is the "triangle inequality" $| x+y \mid \leq$ $|x|+|y|$, so named because in the case of the complex numbers C it says that in the complex plane one side of a triang1e is less than or equal to the sum of the other two sides. The norm $\left|\left.\right|_{p}\right.$ on $Q_{p}$ satisfies a stronger inequality:

$$
|x+y|_{p}^{p} \leq \max \left(|x|_{p},|y|_{p}\right)
$$

This is obvious if we recall how to evaluate $|x|_{p}$ for $x=$ $a_{m} p^{m}+a_{m+1} p^{m+1}+\ldots$ (see above). A norm that satisfies (2.1) is called "non-Archimedean". Inequality (2.1) is sometimes called the "isosceles triangle principle", because it immediately implies that, among the three "sides" $|x| p,|y|_{p}$ and $|x+y|_{p}$; at least two must be equal. Thus, in non-Archimedean geometry "all triangles are isosceles".

Here is another strange consequence of (2.1). In a field with a non-Archimedean norm $\left|\left.\right|_{p}\right.$, define

$$
\begin{aligned}
& D_{a}(r)=\left\{x| | x-\left.a\right|_{p} \leq r\right\} \quad(\text { "closed" disc of radius } \\
& r \text { centered at a) }
\end{aligned}
$$

Then if $b \in D_{a}(r)$, it follows from (2.1) that $D_{b}(r)=D_{a}(r)$.
(Also, if $b \in D_{a}\left(r^{-}\right)$, then $D_{b}\left(r^{-}\right)=D_{a}\left(r^{-}\right)$.) Thus, any point in a disc is its center! In particular, any point in a disc (or in its complement) has a neighborhood completely contained in the disc (resp., in its complement). Therefore, any disc is both open and closed in the topological sense. That is why the words "open" and "ciosed" In (2.2) are in quotation marks; these words are used only by analogy with classical geometry, and one should not be misled by them.

In $p_{p}$, it is not hard to see that all discs of finite radius are compact. The most important such disc is
$Z_{p \text { déf }} D_{0}(1)=\left\{\left.x| | x\right|_{p} \leq 1\right\}=\left\{x=a_{0}+a_{1} p+a_{2} p^{2}+\ldots\right\}$.
$Z_{p}$ is a ring, whose elements are called p-adic "integers". $Z_{p}$ is the closure of the ordinary integers $Z$ in $Q_{p}$. In $Q_{p}$, the other discs centered at 0 are

$$
p^{m} z_{p}=\left\{x=a_{m} p^{m}+a_{m+1} p^{m+1}+\ldots\right\} \quad \text { for } m \in z
$$

$Z_{p}$ is a local ring, i.e., it has a unique maximal ideal $p Z_{p}$, and its residue field $Z_{p} / p Z_{p}$ is the field of $p$ elements $F_{p}=z / p Z$. The set of invertible elements in the ring $Z_{p}$ is
$Z_{p}^{*} \underset{d e f}{=} Z_{p}-P_{p}=\left\{\left.x| | x\right|_{p}=1\right\}$

$$
=\left\{x=a_{0}+a_{1} p+a_{2} p^{2}+\ldots \mid a_{0} \neq 0\right\} .
$$

There are $p-1$ numbers in $\underset{p}{Z *}$ which play a special role: the ( $p-1$ )-th roots of one. For each possible choice of $a_{0}=$ $1,2, \ldots, p-1$, there is a unique such root whose first digit is ${ }^{a} a_{0}$; we denote it $\omega\left(a_{0}\right)$ and call it the Teichmüler representative of $a_{0}$. For example, for $p=5$

$$
\begin{aligned}
& \omega(1)=1 \\
& \omega(2)=2+1 \cdot 5+2 \cdot 5^{2}+1 \cdot 5^{3}+3 \cdot 5^{4}+\ldots \\
& \omega(3)=3+3 \cdot 5+2 \cdot 5^{2}+3 \cdot 5^{3}+1 \cdot 5^{4}+\ldots=-\omega(2) \\
& \omega(4)=4+4 \cdot 5+4 \cdot 5^{2}+4 \cdot 5^{3}+4 \cdot 5^{4}+\ldots=-1 .
\end{aligned}
$$

Except for $\omega( \pm 1)$, the Teichmuller representatives are irrational, so their p-adic digits do not repeat, and can be expected to be just as random as, say, the decimal digits in $\sqrt{2}$.

$$
\text { If } x=a_{0}+a_{1} p+\ldots \in Z_{p}^{*} \text {, we set } \omega(x)=\omega\left(a_{0}\right) \text {. Any }
$$ write

$$
x=p^{\operatorname{ord}_{p} x} w\left(x_{0}\right)<x_{0}>
$$

where $\left\langle\mathrm{x}_{0}\right\rangle \mathrm{d}=\mathrm{ef} \mathrm{x}_{0} / \omega\left(\mathrm{x}_{0}\right)$ is in $1+\mathrm{p}_{\mathrm{p}}$, the set of x such that $|x-1|_{p}<1$.

The ring $Z_{p}$ is the inverse limit of the rings $Z / p^{n} Z$ with respect to the map "reduction mod $p^{n_{11}}$ from $Z / p^{m_{Z}}$ to $Z / p^{n_{Z}}$ for $\mathrm{m} \geq \mathrm{n}$. This suggests that, if we want to solve an equation $\mathrm{f}(\mathrm{x})=$ 0 for $x \in Z_{p}$, we should first solve it in $Z / p Z=F_{p}$, then in $z / p^{2} Z, z / \mathrm{p}^{3} \mathrm{Z}$, and so on. An important condition under which a solution in $F_{p}$ can be "lifted" to a solution in $Z_{p}$ is given by

Hensel's Lemma. Suppose that $f(x) \in Z_{p}[x], f\left(a_{0}\right) \equiv 0(\bmod p)$, and $f^{\prime}\left(a_{0}\right) \not \equiv 0(\bmod p)$ (here $f^{\prime}$ is the formal derivative of the polynomial f). Then there exists a unique $x=a_{0}+\ldots \in Z_{p}$ such that $f(x)=0$

Hensel's Lemma is proved by Newton's method for approximating roots (see $[59,53]$ ).

$$
\text { For example, when } f(x)=x^{p-1}-1, \quad \text { any } a_{0} \in\{1, \ldots, p-1\}
$$ satisfies $f\left(a_{0}\right) \equiv 0(\bmod p)$, while $f^{\prime}\left(a_{0}\right)=(p-1) a_{0}^{p-2} \not \equiv 0$ (mod p); so Hense1's Lemma tells us that $a_{0}$ has a unique Teichmïller representative $\omega\left(a_{0}\right) \in Z_{p}^{*}$.

Unlike in the case of $R$, whose algebraic closure $C$ is only a quadratic extension, $Q_{p}$ has algebraic extensions of arbitrary
degree; its algebraic closure $\bar{Q}_{p}$ has infinite degree over $Q_{p}$ Can $\left|\left.\right|_{p}\right.$ be extended from $Q_{p}$ to $\bar{Q}_{p}$ ? We11, suppose $\alpha$ is ${ }^{p}$ algebraic over $Q_{p}$ and satisfies the minimal polynomial $f(x)=$ $x^{d}+a_{d-1} x^{d-1}+\cdots+a_{0} . \quad$ It is not hard to show that a multiplicative norm on $\bar{Q}_{p}$ extending $\left|\left.\right|_{p}\right.$ would have to be unique. So the value of this extended $\left|\left.\right|_{p}\right.$ on $\alpha$ and each of its conjugates would be the same (because we can also get an extension of $\left.\left|\left.\right|_{p}\right.$ by composing our first extension of $|\right|_{p}$ with a field automorphism of $\overline{\mathrm{Q}}_{\mathrm{p}}$ taking $\alpha$ to the conjugate). Therefore, the only possible value for $|\alpha|_{p}$ is the $d$-th root of $\left|a_{0}\right|_{p}$. It turns out that this definition

$$
|\alpha|_{p}=\sqrt[d]{\left|N_{Q_{p}}(\alpha) / Q_{p}(\alpha)\right|_{p}} \quad(N \text { denotes field norm) }
$$

does in fact give a norm on $\bar{Q}_{p}$. But the fact that this $\left|\left.\right|_{p}\right.$ satisfies the triangle inequality is not trivial to prove. The extension of $\mid l_{p}$ to $\bar{Q}_{p}$ is perhaps the hardest of the basic facts about p-adic numbers; for two different proofs, see [13] and [53].

We now define the ord function on $\bar{Q}_{p}$ by ord $\alpha=-\log _{p}|\alpha|_{p}$,
as to agree with the earlier ord on $Q_{p}$. (Here $\log$ is the so as to agree with the earlier ord $p_{p}$ on $Q_{p}$. (Here $\log _{p}$ is the orsinary "log to base p ", not to be confused with a p-adic logarithm which we shall introduce shortly.) Clearly, if $\left[\mathrm{K}: \mathrm{Q}_{\mathrm{p}}\right]=$ $d$, then the image of $K$ under ord ${ }_{p}$ is an additive subgroup of $\frac{1}{d} z$, and so ord ${ }_{p} K=\frac{1}{e} Z$ for some $e^{p}$ dividing $d$. This positive integer $e$ is called the index of ramification of $K$. There are two extremes:
(1) $e=1$. Then $K$ is called unramified. An example is $K=Q_{p}(\sqrt[N]{1})$ for $N$ not divisible by $p$. In fact, it can be shown that every unramified $K$ is contained in some cyclotomic field, so the "unramified closure" of $Q_{p}$ is $Q_{p}^{\text {unr }}=\bigcup_{p / N} Q_{p}(\sqrt{1})$.
(2) $e=d$. Then $K$ is called totally ramified. An example is $K=Q_{p}(\xi)$ for $\xi \neq 1$ a $p$-th root of one, i.e., a root of $x^{P^{-1}}+x^{p}{ }^{-2}+\ldots+x+1=0$. To show that $K$ is totally rami-
fied, it suffices to find $\lambda \in K$ such that ord ${ }_{p}=1 /(p-1)$. Let $\lambda=\xi-1$. Since $\lambda$ satisfies: $0=\left[(x+1)^{p}-1\right] /[(x+1)-1]=x^{p-1}+$ $p x^{p-2}+\frac{1}{2} p(p-1) x^{p-3}+\ldots+p$, it follows that ord $p^{\lambda}=\frac{1}{p-1}$ ord $_{p} p$ $=1 /(p-1)$. More generally, if $\xi$ is a primitive $p^{n}$-th root of one, then $Q_{p}(\xi)$ is totally ramified of degree $p^{n}-p^{n-1}$, and

$$
\operatorname{ord}_{p}(\xi-1)=\frac{1}{p^{n}-p^{n-1}}
$$

The set of all totally ramified extensions is harder to describe than the set of all unramified extensions. And, of course, "most" extensions are neither unramified nor totally ramified. In the general case we write $d=e \cdot f$.

The significance of $£$ is as follows. If $K$ is any field with a non-Archimedean norm $\left|\left.\right|_{p}\right.$, we let

$$
0_{K}=\left\{\left.x \in K|\quad| x\right|_{p} \leq 1\right\}, \quad M_{K}=\left\{\left.x \in K| | x\right|_{P}<1\right\}
$$

$0_{K}$ is called the "ring of integers" of $K$, and $M_{K}$ is the unique maximal ideal in $O_{K}$. If $K$ is algebraic over $Q_{P}$, then the residue field $0_{K} / M_{K}$ will be algebraic over $F_{p}$. If $K$ has degree $d$ and ramification index $e$, then this residue field has degree $f=$ d/e over $F_{p}$ (see [59]).

Let us return to the case of $K$ unramified, of degree $d=f$. Let $q=p^{f}$, so that $0_{K} / M_{K}$ is the field of $q$ elements $F_{q}$. Then, using Hensel's Lemina (generalized to $O_{K}$ ), we see that every nonzero element $a_{0} \in \mathrm{~F}_{\mathrm{q}}$ has a unique Teichmüller representative $\omega\left(a_{0}\right) \in K$ such that $\omega\left(a_{0}\right)^{q-1}=1$ and $\omega\left(a_{0}\right) \bmod M_{K}$ is $a_{0}$. If $a_{0}$ generates $F_{q}$ as an extension of $F_{p}$, then $K=Q_{p}\left(\omega\left(a_{0}\right)\right)$. These Teichmüller representatives are a natural choice of "digits" in $K$ : every $\mathrm{x} \in \mathrm{K}$ can be written uniquely as the limit of a sum
$x=\sum_{i>m} a_{i} p^{i}$, where $a_{i} \in\{\omega(a)\}_{a \in F_{q}}$
(we agree to let $\omega(0)=0$ ). Even in $Q_{p}$ it is sometimes convenient to choose $0, \omega(1), \omega(2), \ldots, \omega(p-1)$ as digits instead of $0,1,2, \ldots, p-1$.

Since the complex number field $\mathcal{C}$ is a finite dimensional $R$-vector space, it is complete under the extension of $\mid$ | to $C$. However, $\bar{Q}_{p}$ turns out not to be complete under $\left|\left.\right|_{p}\right.$. For example, the convergent infinite sum $\Sigma \mathrm{x}_{\mathrm{i}} \mathrm{p}^{1}$, where the $\mathrm{x}_{\mathrm{i}}$ are a sequence of roots of one of increasing degree, in general is not algebraic over $Q_{p}$. Thus, in order to do analysis, we must take a larger field than $\bar{Q}_{p}$. We denote the completion of $\bar{Q}_{p}$ by $\Omega_{p}$ : $\Omega_{\mathrm{p}}=\hat{\bar{Q}}_{\mathrm{p}} \quad\left(\wedge\right.$ means completion with respect to $\left.\mid \|_{\mathrm{p}}\right)$.
It is not hard to see that $\Omega_{\mathrm{p}}$ is algebraically closed, as well as complete, that $0_{\Omega_{p}} / M_{\Omega_{p}}=\bar{F}_{p}$, and that ord $\Omega_{p}=Q$. Sometimes $\Omega_{p}$ is denoted $C_{p}$ in order to emphasize the analogy with the complex numbers (i.e., both are the smallest extension field of $Q$ that is both algebraically closed and complete in the respective metric). But in some respects $\Omega_{p}$ is more complicated. For example, it is a much bigger extension of $Q_{p}$ than $C$ is over $R$ (in fact, $\Omega_{p}$ has uncountable transcendance degree over $Q_{p}$ ), and it is easy to see that $\Omega_{p}$ is not locally compact.
3. Power series

An infinite sum $\Sigma a_{i}$ has a limit if $\Sigma_{N \leq i<M} a_{i}$ is small for large $N, M>N$. Because of the isosceles triangle principle (2.1), in $\Omega_{p}$ this occurs if and only if $a_{i} \rightarrow 0$, i.e., $\left|a_{i}\right|_{p}$ $\rightarrow 0$, or equivalently, ord ${ }_{p}{ }_{i} \rightarrow \infty$. Thus, the question of convergence or divergence of a power series $\sum a_{i} x^{i}$ depends only on $|x|_{p}$, not on the precise value of $x$. There is no "conditional convergence". Thus, every infinite series $\sum a_{i} x^{i}$ has a radius of convergence $r$ such that one of the following holds:
$\sum_{i=0}^{\infty} a_{i} x^{i}$ converges $\Leftrightarrow x \in D\left(r^{-}\right) \quad\left({ }_{d e f} D_{0}\left(r^{-}\right)\right.$, see (2.2))
or

$$
\left.\sum_{i=0}^{\infty} a_{i} x^{i} \text { converges } \Leftrightarrow x \in D(r) \quad \zeta_{d e f} D_{0}(r)\right) .
$$

An example of the first alternative is $\Sigma \mathrm{x}^{\mathrm{i}}$ (where $\mathrm{r}=1$ ); an
example of the second is the derivative $\Sigma \mathrm{p}^{i} \mathrm{x}^{\mathrm{p}^{i}-1}$ of $\Sigma \mathrm{x}^{\mathrm{p}^{i}}$ (here also $r=1$ ).

An important example is the series $e^{x}=\Sigma x^{i} / i!$. To determine its radius of convergence, we must find ord ${ }_{p}(i)$. If $i$ is a power of $p$, it is easy to see that ord $\left(p^{n}!\right)=p^{n-1}+p^{n-2}+$ $\ldots+p+1=(i-1) /(p-1)$. More generally, if we write the positive integer $i$ to the base $p: i=\Sigma a_{i} p^{i}$, and let $s_{i}=\Sigma a_{i}$ denote the sum of its digits, then

$$
\begin{equation*}
\operatorname{ord}_{p}(i!)=\frac{i-S_{i}}{p-1} \tag{3.1}
\end{equation*}
$$

Since $1 \leq s_{i} \leq(p-1)\left(\log _{p} i+1\right)$, it follows that asymptotically $\operatorname{ord}_{p}(i!) \sim \frac{i}{p-I}, \quad$ and so

$$
\begin{aligned}
\operatorname{ord}_{p}\left(x^{i} / i!\right) \longrightarrow \infty & \Leftrightarrow \operatorname{ord}_{p} x>\frac{1}{p-1} \\
& \Leftrightarrow x \in D\left(\frac{1}{\gamma}\right), \quad \gamma=\frac{p-1}{p}>1
\end{aligned}
$$

Thus, $e^{x}$ converges in a disc smaller than the unit disc. In the classical case the i! in the denominator makes $e^{x}$ converge everywhere, but in $\left|\left.\right|_{p}\right.$ it has a harmful effect on convergence. The poor convergence of $e^{x}$ causes much of p-adic analysis, e.g., differential equations, to involve subtleties which are absent in complex analysis.

To obtain a series convergent in $D\left(1^{-}\right)$instead of $D\left(\frac{1-}{\gamma}\right)$, we can replace $e^{x}$ by $e^{\pi x}$, where $\pi$ (not to be confused with the real number $\pi=3.14 \ldots$ ) is any element of $\Omega_{p}$ such that ord $\pi$ $=1 /(p-1)$. The best choice of $\pi$ is a ( $p-1)-$ th root of $-p$, for reasons that will become clear later.

We can analyze more closely why $e^{x}$ converges so poorly if we use the formal power series identity

$$
\begin{equation*}
e^{x}=\prod_{n=1}^{\infty}\left(1-x^{n}\right)^{-\mu(n) / n} \quad \text { in } Q[[x]] \tag{3.2}
\end{equation*}
$$

where the Möbius function $\mu$ is defined by
$\mu(\mathrm{n})=\left\{\begin{array}{l}0 \quad \text { if there is a prime whose square divides } n ; \\ (-1)^{k} \text { if } n \text { is a product of } k \text { distinct primes. }\end{array}\right.$
The identity (3.2) is easily proved by taking log of both sides and using the fact that $\sum_{d \mid n} \mu(d)=1$ if $n=1$ and 0 otherwise.

Most of the terms in (3.2) -- those for which p/n -- have fairly good convergence, because the binomial series

$$
(1+Y)^{\alpha}=\sum\binom{\alpha}{i} Y^{i}, \quad\binom{\alpha}{i}=\frac{\alpha(\alpha-1) \cdots(\alpha-i+1)}{i!}
$$

has coefficients

$$
\begin{equation*}
\binom{\alpha}{i} \in Z_{p} \quad \text { for } \quad \alpha \in Z_{p} \tag{3.3}
\end{equation*}
$$

(Namely, this is trivial for $\alpha$ a positive integer; then use the fact that the positive integers are dense in $z_{p}$.) Thus, for $\mathrm{p} / \mathrm{n}, \prod^{\boldsymbol{T}} \mathrm{I}^{\left(1-\mathrm{x}^{\mathrm{n}}\right)^{-\mu(n) / n} \in Z_{p}[[\mathrm{x}]]}$ and so converges for $|\mathrm{x}|_{\mathrm{p}}<1$. The bad convergence of (3.2) comes from those $n$ which are divisible by $p$. So, to get better convergence, we can define the "Artin-Hasse exponential"

$$
E_{p}(x)=\prod_{P} / n=x^{n}\left(1-x^{n}\right)^{-\mu(n) / n}=e^{x+x^{p} / p+x^{p} / p^{2}+x^{p} / p^{3}+\ldots}
$$

where the last equality of formal power series is proved in the same way as (3.2). Then $E_{p}(x)$ is in $Z_{p}[[x]]$, and so converges in $D\left(1^{-}\right)$.

If we make the change of variables $E_{p}(\pi x)$, where $\pi^{p-1}=-p$, the first two terms in the exponent are $\pi\left(x-x^{p}\right)$. The expression $x-x^{p}$ plays a key role in much of p-adic analysis, since in a field of characteristic $p$

$$
\mathrm{x} \in \text { the prime field } \mathrm{F}_{\mathrm{p}} \Leftrightarrow \mathrm{x}-\mathrm{x}^{\mathrm{p}}=0
$$

also recall that the Teichmüller representatives $\{\omega(a)\}_{a \in F_{p}}$ are solutions of this equation in $Q_{p} . \quad \pi$ is chosen to be a $(p-1)$-th root of $-p$ precisely so that the first two terms in the exponent for $E_{p}(\pi x)$ become a multiple of $x-x^{p}$.

$$
\begin{align*}
& \text { Since } \\
& e^{\pi\left(x-x^{p}\right)}=E_{p}(\pi x) \prod_{\pi \geq 2} e^{-(\pi x)^{p^{i}} / p^{i}},
\end{align*}
$$

the convergence of $e^{\pi\left(x-x^{P}\right)}$ is determined by the worst convergence that occurs on the right. $E_{p}(\pi x)$ converges on $D\left(\gamma^{-}\right)$(recall $\left.Y=p^{1 /(p-1)}\right)$, and it is easy to compute that the worst series is the first one in the product, $\exp \left(-\pi^{2} x^{p} p^{2} / p^{2}\right)$, which converges for ord ${ }_{\mathrm{p}} \mathrm{x}>-(\mathrm{p}-1) / \mathrm{p}^{2}$. Thus, if we let $\gamma_{1}=\mathrm{p}^{(\mathrm{p}-1) / \mathrm{p}^{2}}>1$, it follows that $e^{\pi\left(x-x^{P}\right)}$ converges on $D\left(\gamma_{1}^{-}\right)$, a disc strictly bigger than $D(1)$

Thus, the $-x^{p}$ in $e^{\pi\left(x-x^{p}\right)}$ is a "correction" which improves the convergence of $e^{\pi x}$. We can see how this works if we look at the expansions of $e^{\pi x}$ and $e^{\pi\left(x-x^{p}\right)}$ out to the $x^{p}$-term, the first term where the two series differ. In the expansion $e^{\pi x}=$ $\sum(\pi x)^{\frac{1}{i}} / i!$, the $x^{P}$-term is the first one in which $1 / i!\& Z_{p}$, i.e., the first term containing a $p$ in the denominator. Thus, $\left|\pi^{p} / p!\right|_{p}=|-\pi /(p-1)!|_{p}=|\pi|_{p}$ (the first equality because $\pi^{p-1}$ $=-p$ ). But the coefficient of $x^{p}$ in $e^{\pi\left(x-x^{p}\right)}$ is

$$
\frac{\pi^{p}}{p!}-\pi=\pi(-p / p!-1)=\frac{-\pi}{(p-1)!}(1+(p-1)!)
$$

A simple fact of elementary number theory (Wilson's theorem) says: $(p-1)!\equiv-1$ (mod $p)$. Hence, the p-adic norm of the coefficient of $\mathrm{x}^{\mathrm{P}}$ in $e^{\pi\left(x-x^{P}\right)}$ is bounded by $|\mathrm{p} \pi|_{p}=\left|\pi^{p}\right|_{p}$. Thus, the correction term $-\pi x^{P}$ has the effect of canceling the $p$ in the denominator of $(\pi x)^{p} / p$.

We denote $E_{\pi}(x)=e^{\pi\left(x-x^{p}\right)}$ (not to be confused with $E_{p}(x)$ ). Note that $E_{\pi}(x)$ must first be expanded as a power series and then evaluated. If $|x|_{p}<1$, the result will be the same as if we first substituted x in $\pi\left(\mathrm{x}-\mathrm{x}^{\mathrm{P}}\right)$ and then took the exponential. But if $|x|_{p} \geq 1$, that exponential will not converge unless $\left|x-x^{P}\right|_{p}<1$, and even in the latter case will in general give the wrong value; for example, $E_{\pi}(1) \neq 1=e^{0}$ (see §III.5).

Another important series is

$$
\begin{equation*}
\log (1+x)=\sum_{i=1}^{\infty} \frac{(-1)^{i+1}}{i} x^{i}, \tag{3.5}
\end{equation*}
$$

which is easily seen to converge on $D\left(1^{-}\right)$. It has better convergence than $e^{x_{!}}$Since the identity
$\log (x y)=\log x+\log y$
holds as a formal power series identity, i.e., $\sum(-1)^{i+1} \mathrm{x}^{\mathrm{i}} / \mathrm{i}+$ $\sum(\sim 1)^{i+1} y^{i} / i=\sum(-1)^{i+1}(x+y+x y)^{i} / i$ in $Q[[x, y]]$, it follows that (3.6) holds in $\Omega_{p}$ as long as $|x-1|_{p}<1$ and $|y-1|_{p}<1$. In particular, since $|\xi-1|_{p}<1$ for $\xi$ any $p^{\text {n }}$-th root of one (see §2), we can apply (3.6) to conclude that $\log \xi=0$.

The p-adic logarithm has a natural extension to $\Omega_{p}^{*}=\Omega_{p}-\{0\}$, which we shall denote $\ln _{p}$ (so as not to confuse it with the classical log-to-the-base-p; however, in the literature $\log _{p}$ is normally used rather than $\ln _{p}$ ).

Proposition. There exists a unique function $n_{p}: \Omega_{\mathrm{p}}^{*} \longrightarrow \Omega_{\mathrm{p}}$ such that
(1) $1 n_{p}(1+x)$ is given by the series (3.5) if $|x|_{p}<1$;
(2) (3.6) holds for all $x, y \in \Omega_{\mathrm{p}}^{*}$;
(3) $\ln _{\mathrm{p}}(\mathrm{p})=0$.

The third condition is a normalization, which is necessary because, as mentioned before, ord ${ }_{p}$ behaves like a logarithm. Thus, if $1 n_{p}$ is any function satisfying (1) and (2), then for any constant $c \in \Omega_{p}$ the function $\ln _{p}+c \cdot o r d_{p}$ also satisfies (1) and (2).

I won't prove this proposition, but will discuss concretely how one computes a logarithm. First, for every $\frac{m}{n} \in Q$, choose ${ }^{p} p^{m / n}$, to be any root of $x^{n}-p^{I I}=0$. Now suppose we want to find $\ln _{p} x$ for some nonzero $x \in \Omega_{p}$. First write $x=p^{m / n} x_{0}$, where $m / n=$ ord $_{p} x$. Since $\left|x_{0}\right|_{p}=1$, its reduction modulo $M_{\Omega_{p}}$ is a nonzero element $\bar{x}_{0} \in \bar{F}_{p}$. Let $\omega\left(\bar{x}_{0}\right)$ be the Teichmüller representative of $\bar{x}_{0}$. Then

$$
\mathrm{x}=\mathrm{p}^{\text {ord }_{\mathrm{p}} \mathrm{x}} w\left(\overline{\mathrm{x}}_{0}\right)<\mathrm{x}_{0}>\text {, where }\left|<\mathrm{x}_{0}\right\rangle-\left.1\right|_{\mathrm{p}}<1 \text {. }
$$

Since $1_{n_{p}} p=0$, and (3.6) implies that $\ln _{p}$ (any root of 1 ) $=0$, we have

$$
\left.\ln _{p} x=1 n_{p}<x_{0}\right\rangle=\sum(-1)^{i+1}\left(\left\langle x_{0}\right\rangle-1\right)^{i} / i
$$

$$
\begin{aligned}
\text { For example, } & \text { ( } \begin{aligned}
& \text { p } \\
& \ln _{5}\left(\frac{1}{250}\right)=\ln _{5}\left(\frac{2+1 \cdot 5+2 \cdot 5^{2}+1 \cdot 5^{3}+3 \cdot 5^{4}+\ldots}{2}\right) \\
&=\sum(-1)^{i+1}\left(3 \cdot 5+3 \cdot 5^{2}+0 \cdot 5^{3}+4 \cdot 5^{4}+\ldots\right)^{i} / i
\end{aligned},
\end{aligned}
$$

Note that a function such as ord ${ }_{p}$, which is locally constant on $\Omega_{p}^{*}$ (i.e., for every $a \in \Omega_{p}^{*}$ there exists $r$ such that ord ${ }_{p} x$ $=$ ord $_{p}$ for $x \in D_{a}(r)$ ) but is not constant, could not exist on $\mathrm{C} *$. For this reason, the theory of analytic continuation is more complicated on $\Omega_{\mathrm{p}}^{*}$. Unlike the classical $\log , \ln _{p}$ is not obtained by "analytic continuation" of the series (3.5); any of the functions $1 n_{p}+c \cdot o r d_{p}$ would also be locally analytic and agree with (3.5) on $\mathrm{D}_{1}\left(1^{-}\right)$.

There is a notion of p-adic global analyticity, due to Krasner [57], such that two globally analytic functions which agree, say, on a disc, must agree everywhere. Nameiy, let $D \subset \Omega_{p}$ be a socalled "quasi-connected" set, the most important examples of which are discs from which finitely many smaller discs and/or compact subsets have been removed. Then a function $f: D \rightarrow \Omega_{P}$ is said to be Krasner analytic if $D$ is a union of open sets $D_{i}, D_{i} \subset D_{i+1}$, such that for each $i,\left.f\right|_{D_{i}}$ is a uniform limit of rational functions having no poles in $D_{i}$. $\ln _{p}$ is not Krasner analytic on $\Omega_{\mathrm{p}}^{*}$. Later we shall see examples of interesting Krasner analytic functions. For example, the second derivative of the p-adic log gamma function turns out to be Krasner analytic on the complement of $Z_{p}: \quad D=\Omega_{p}-Z_{p}$ (see p. 134).

A final remark about $\ln _{p}$ : it has the expected derivative $\frac{1}{x}$, since $\lim _{\varepsilon \rightarrow 0}\left[\left(\ln _{p}(x+\varepsilon)-\ln _{p} x\right) / \varepsilon\right]=\lim \frac{1}{\varepsilon} \ln _{p}\left(1+\frac{\varepsilon}{x}\right)$, and $\ln _{\mathrm{p}}\left(1+\frac{\varepsilon}{\mathrm{x}}\right.$ ) is given by the usual series as soon as $|\varepsilon|_{\mathrm{p}}<|\mathrm{x}|_{\mathrm{p}}$.
4. Newton polygons
a. Classical case

For $f(X, Y)=\sum a_{i j} X^{i} Y^{j} \in R[X, Y]$, let $M_{f}$ be the convex hull
the following set of points in the (i, $j)-p l a n e:\left\{(i, j) \mid a_{i j} \neq 0\right\}$. of the following set of points in the $(i, j)$-plane: $\left\{(i, j) \mid a_{i j} \neq 0\right\}$. $M_{f}$ is called the Newton polygon of $f$. If two polynomials $f, g$ $\epsilon \mathrm{R}[\mathrm{X}, \mathrm{Y}]$ have no common factors, then the two curves determined by $f$ and $g$ intersect in a finite number $N$ of points (counting multiplicity $):\{(x, y) \mid f(x, y)=g(x, y)=0\}$. Let $\quad M_{f}+M_{g}=$ $\left\{z=x+y \mid x \in M_{f}, y \in M_{g}\right\}$. Then it can be shown that
$\mathrm{N} \leq \operatorname{area}\left(\mathrm{M}_{\mathrm{f}}+\mathrm{M}_{\mathrm{g}}\right)-\operatorname{area}\left(\mathrm{M}_{\mathrm{f}}\right)-\operatorname{area}\left(\mathrm{M}_{\mathrm{g}}\right)$.
b. The p-adic case: polynomials

Let $f(x)=a_{0}+\ldots+a_{d} x^{d} \in \Omega_{p}[x]$. The Newton polygon $M_{f}$ of $f$ is defined to be the convex hull of the points ( $i$, ord $p_{i} a_{i}$ ) (where we agree to take ord ${ }_{p} 0=+\infty$ ), i.e., $M_{f}$ is the polygonal line obtained by rotating a vertical line through $\left(0\right.$, ord $\left._{p} a_{0}\right)$ counterclockwise until it bends around various points (i, ord $a_{i}$ ), and eventually reaches the point
 (d, ord ${ }_{p} a_{d}$ ). This is similar to the classical case (where we take $a_{i}=\stackrel{\sum}{j}_{j} a_{i j} Y^{j}$ and $\operatorname{ord}_{Y} a_{i}=$ the least $j$ for which $a_{i j} \neq 0$ ), except that we only take the lower part of the convex hull.

It is not hard to prove the following
Proposition. If a segment of ${ }^{M} f$ has slope $\lambda$ and horizontal length $N$ (i.e., it extends from ( $i$, ord $p_{i}$ ) to $\left(i+N, \lambda N+\operatorname{ord}_{p} a_{i}\right)$, then $f$ has precisely $N$ roots $r_{i}$ with ord ${ }_{p} r_{i}=-\lambda \quad$ (counting multiplicity).

Examples. (1) The Efsenstein irreducibility criterion: if $f(x)=a_{0}+\ldots+a_{d-1} x^{d-1}+x^{d} \in Q[x]$, and if there exists a prime $p$ such that $\operatorname{ord}_{p} a_{i} \geq 1$ for $0 \leq i<d$ and ord $a_{0}=1$, then $f$ is irreducible over $Q$. In fact, using the Newton polygon $M_{f}$, we can quickly see that $f$ is even irreducible over $Q_{p}$. Namely, the
conditions on ord ${ }_{p}{ }_{i}$ imply that $M_{f}$ consists of the line segment from $(0,1)$ to $(d, 0)$. Hence $£$ has $d$ roots all of ordinal $\frac{1}{d}$. If $f$ factored over $Q_{p}$, each root $r$ would have degree $d^{+}<d$ over $Q_{p}$, and hence we would have ord $p \in \frac{1}{d} r^{2}$. Thus, $f$ is irreducible.
(2) Later we'11 want to study the curve $y^{p}-y=x^{d}$. If this curve is considered over a field of characteristic $p$, there are $p$ obvious automorphisms $\mathrm{x} \mapsto \mathrm{x}, \quad \mathrm{y} \mapsto \mathrm{y}+\bar{a}, \quad \overline{\mathrm{a}} \in \mathrm{F}_{\mathrm{p}} . \quad$ Suppose we want to find similar automorphisms $\mathrm{x} \mapsto \mathrm{x}, \mathrm{y} \mapsto \mathrm{y}+\mathrm{a}$ when the curve is considered over $\Omega_{p}$. For example, let us fix $y \in \Omega_{p}$ and look for $a \in \Omega_{p}$ such that sending $y \mapsto y+a$ "lifts" the automorphism $y \mapsto y+1$ in the sense that $a=1+z$ with $|z|_{p}<1$, i.e., $a \equiv 1\left(\bmod M_{\Omega}\right)$. It is convenient to suppose that $|y|_{p}<\gamma$, where $\gamma=p^{1 /(p-1)} \stackrel{p}{>}$. We must choose $z$ so that

$$
(y+1+z)^{p}-(y+1+z)=y^{p}-y
$$

or, if we write this as a polynomial in $z$,

$$
z^{p}+\sum \underset{i=1}{p-2}\binom{\mathrm{p}}{\mathbf{i}}(\mathrm{y}+1)^{\mathrm{i}} z_{z^{p-i}}+\left(\mathrm{p}(\mathrm{y}+1)^{\mathrm{p}-\mathrm{i}}-1\right) z+
$$

$$
+\left[(y+1)^{p}-y^{p}-1\right]=0
$$

The constant term $a_{0}=(y+1)^{p}-y^{p}-1=\Sigma_{1 \leq i<p}\binom{\mathrm{p}}{i} y^{i}$ satisfies $\operatorname{ord}_{p} a_{0} \geq 1+\min \left(0,(p-1) \operatorname{ord}_{p} y\right)$, which is greater than zero, since we have assumed that ord ${ }_{p} y>-1 /(p-1)$. On the other hand, ord $a_{1}$ $=\operatorname{ord}_{p} a_{p}=0$ and ord $a_{i} \geq 0$ for $1<i<p$. Hence, the Newton polygon of this polynomial in $z$ is as shown in the diagram to the right. The only nonzero slope is the first little seg-
 ment, with slope $\lambda=-\operatorname{ord}_{p} a_{0}$. Thus, there is exactiy one root $z$ with $|z|_{p}<1$, in fact, with ord $p_{p}=-\lambda=$ ord $a_{0}$. This root $z$ gives the unique Iifting to $\Omega p$ of the automorphism $y \mapsto y+1$ in characteristic $p$. The other $p-1$ roots $z$ have $|z|_{p}=1$, and the corresponding maps $y \mapsto y+1+z$ lift the other automorphisms $y \mapsto y+\bar{a}, \quad \bar{a} \in F_{p}$.
c. The p-adic case: power series

The Newton polygon $M_{f}$ for a power series $f(x)=\sum a_{i} x^{i}$ $\epsilon \Omega_{p}[[x]]$ is defined just as for polynomials, except that now it extends infinitely far to the right. Also, it is possible for the Newton polygon to include an infinitely long segment without any points ( $i$, ord $a_{i}$ ) far to the right. For example, the power series $1+\Sigma_{j>1} p^{j-1} x^{p^{j}}$ has simply the $x$-semiaxis as its Newton polygon, although ord $a_{i}>0$ for $i>p$. Here is the case $p=2$ :


The following theorem is the p-adic analog of the Weierstrass Preparation Theorem.

Theorem. Let $f(x)=a_{m} x^{m}+\ldots \in \Omega_{p}[[x]], a_{m} \neq 0$, be a power series which converges on $D\left(p^{\lambda}\right)$. Let ( $N$, ord ${ }_{p} a_{N}$ ) be the right endpoint of the last segment of $M_{f}$ with slope $\leq \lambda$, if this $N$ is finite. Otherwise, there will be a last infinitely long segment of slope $\lambda$ and only finitely many points ( $i$, ord ${ }_{p} a_{i}$ ) on that segment. In that case let $N$ be the last such $i$ (for example, in the above illustration $N=2$ ). Then there exists a unique polynomial $h(x)$ of the form $b_{m} x^{m}+b_{m+1} x^{m+1}+\ldots+b_{N} x^{N}$ with $b_{m}=a_{m}$ and a unique power series $g(x)$ which converges and is nonizero on $D(p)$, such that

$$
f(x)=\frac{h(x)}{g(x)} \text { on } D\left(p^{\lambda}\right) .
$$

In addition, $M_{h}$ coincides with $M_{f}$ as far as the point ( $\mathrm{N}, \operatorname{ord}_{\mathrm{p}}^{\mathrm{a}} \mathrm{a}_{\mathrm{N}}$ ).

Corollary 1. Within the region of convergence of $f$, the Newton polygon determines ord $p$ of the zeros of $f$ in the same way

## as for polynomials.

Corollary 2. A power series which converges everywhere and has no zeros is a constant.

For proofs of these facts, see, for example, [53].
Examples. (1) The power series $1+\Sigma_{j>I} p^{j-1} \mathrm{x}^{j}$, which converges on $D(1)$, has precisely $p$ zeros, all with $\left|\left.\right|_{p}=1\right.$.
(2) For the $\log$ series $f(x)=\sum(-1)^{i+1} x^{i} / i, \quad M_{f}$ is the polygonal line comnecting the points $\left(p^{j},-j\right), j=0,1,2, \ldots$ The picture for $p=2$ is given below. We may conclude that in

$D_{1}\left(1^{-}\right)$the function $1 n_{p}$ vanishes at points $1+x$ for exactly $p^{j}-p^{j-1}$ values of $x$ with ordinal $1 /\left(p^{j}-p^{j-1}\right)$. These $x^{\prime} s$ are precisely $\mathrm{x}=\xi-1$ for $\xi$ a primitive $\mathrm{p}^{j}-\mathrm{th}$ root of one (see (2.3)).

Remark. Some specialists prefer another definition of the Newton polygon. Instead of the points ( $i$, ord $p_{p}$ ), they look at the lines $l_{i}: y=i_{x}+\operatorname{ord}_{p} a_{i}$ with slope $i$ and $y$-intercept ord $a_{i}$. Then $\tilde{M}_{f}$ is defined as the graph of the function $\min _{i} l_{i}(x)$. The $x$-coordinates of the points of intersection of the $\ell_{i}$ which appear in $\tilde{M}_{f}$ give ord ${ }_{p}$ of the zeros, and the difference between the slopes $i$ of successive $l_{i}$ which appear in $\tilde{M}_{f}$ give the number of zeros with given ord ${ }_{p}$. For example, $\tilde{M}_{f}$ for the log series $f(x)=\Sigma(-I)^{i+1} x^{i} / i \quad$ is shown in the drawing on the next page. It somewhat resembles the usual graph of $\log$, especially near the y-axis. This type of Newton polygon was used in Hà-huy-Khoái's thesis [39], which contains a detailed discussion of such Newton polygons, as well as a new generalization of Newton polygons
("Newton sequences") which can be used for more refined investigations of power series.

II. p-ADIC $\zeta$-fUNCTIONS, L-FUNCTIONS, AND $\Gamma$-FUNCTIONS

1. Dirichlet L-series

We leave p-adics for a moment to review the basic facts about Dirichlet L-series (see [13,41]). Let $f: Z \rightarrow C$ be a periodic function with period $d: f(x+d)=f(x)$. Then we define

$$
L(s, f)=\sum_{n=1}^{\infty} f(n) n^{-s}
$$

for $\operatorname{Re} s>1$, and extend by analytic continuation to other $s \in C$.

$$
\begin{align*}
& \text { The generalized Bernoulli numbers are } \\
& \qquad B_{k, f}=k!\cdot \text { coefficient of } t^{k} \text { in } \sum_{a=0}^{d-1} \frac{f(a) t e^{a t}}{d t} .1 \tag{1.1}
\end{align*}
$$

It can be shown [41] that for $k$ a positive integer

$$
\begin{equation*}
L(I-k, f)=-\frac{B_{k, f}}{k} . \tag{1.2}
\end{equation*}
$$

For example, for the Riemann zeta function $\zeta(\mathrm{s})=\mathrm{L}(\mathrm{s}, 1)$ (where 1 denotes the constant function $I$, having period 1 )

$$
\zeta(1-k)=-\frac{1}{k} B_{k}, \quad B_{k}=k!\cdot \text { coefficient of } t^{k} \quad \text { in } \frac{t}{e^{t}-1} .
$$

When $\mathrm{f}=\mathrm{X}$ is a character, i.e., a homomorphism $\mathrm{X}:(\mathrm{Z} / \mathrm{dZ})$ * $\longrightarrow C^{*}$ from the multiplicative group of integers mod $d$ (where $X$ is extended by $X(n)=0$ for all $n$ having a common factor with d), the l-series equals the following "Euler product" if $\operatorname{Re} s>1$ :

$$
\begin{equation*}
L(s, x)=\prod\left(1-\frac{x(l)}{e^{s}}\right)^{-1}, \tag{1.3}
\end{equation*}
$$

where the product is taken over all primes $\ell$.
-functions occur in many situations in number theory. To give
a simple example, the class number $h$ of an imaginary quadratic field $Q(\sqrt{-d})$ of discriminant $-d$ is given by

$$
h=\frac{W \sqrt{d}}{2 \pi} L(1, \chi)=-\frac{w}{2} \frac{1}{d} \sum_{a=1}^{d-1} a \chi(a),
$$

where $w=2,4$, or 6 is the number of roots of unity in $Q(\sqrt{-d})$, and $X:(Z / d Z) * \longrightarrow\{ \pm 1\}$ is the Legendre symbol (quadratic residue symbol). By the way, no elementary proof (not using the Dirichlet formula $h=-\frac{1}{d} \sum$ a $X(a)$ ) is known for the nonvanishing of the simple sum $\Sigma$ a $\chi(a)$. Later (Chapter IV) we shall study generalizations and p-adic analogs of the formula $L(1, X)=$ $2 \pi \mathrm{~h} / \mathrm{w} / \mathrm{d}$.

We shall also want to consider "twisted" L-functions. Let $r$ be a positive integer, and let $\varepsilon \neq 1$ be any nontrivial r-th root of one. Let $z^{d}=\varepsilon$. Then let

$$
\mathrm{L}(\mathrm{~s}, \mathrm{f}, \mathrm{z})=\sum_{\mathrm{n}=1}^{\infty} \mathrm{f}(\mathrm{n}) \mathrm{z}^{\mathrm{n}} \mathrm{n}^{-\mathrm{s}} .
$$

Since the function $n \mapsto f(n) z^{n}$ has period $d r$, this is a special case of the L-series considered above. In particular, if we replace $k$ by $k+1$ and $f(a)$ by $z^{a} f(a)$ in (1.1) and (1.2), we obtain:

$$
\begin{align*}
L(-k, f, z) & =k!\cdot \text { coefficient of } t^{k} \text { in } \sum_{\substack{0 \leq a<d \\
0 \leq b<r}} \frac{f(a) e^{a t} e^{b d t} z a+b d}{1-e^{r d t}} \\
& =k!\cdot \text { coefficient of } t^{k} \text { in } \sum_{0 \leq a<d} \frac{£(a) z^{a} e^{a t}}{1-\varepsilon e^{d t}} . \quad \text { (1.4. }
\end{align*}
$$

We now proceed to the p-adic theory

## 2. p-adic measures

Let $d$ be a fixed positive integer, and let $x=\frac{1 i m}{N} 2 / \mathrm{dp}^{N} Z$, where the map from $Z / d p^{M_{Z}}$ to $Z / d p{ }^{N} Z$ for $M \geq N$ is reduction $\bmod d p{ }^{N}$. In the special case $d=1, X$ is simply $Z_{p}$. By $a+d p Z_{p}$ we mean the set of all $x \in X$ which map to $a$ under the natural map $X \longrightarrow Z / d p^{N} Z$. Without loss of generality, we may agree always to choose $a$ so that $0 \leq a<\mathrm{dp}^{\mathrm{N}}$. Note that
$x=\bigsqcup_{0 \leq a<d} a+d z_{p} \quad$ is a disjoint union of $d$ topological spaces isomorphic to $Z_{\mathrm{P}}$, Also,

$$
a+d p{ }^{N} Z_{p}=\bigsqcup_{0 \leq b<p}\left(a+b d p{ }^{N}\right)+d p^{N+1} Z_{p} \quad \text { (disjoint union). (2.1) }
$$

It is not hard to show that any open subset which is compact
i.e., closed, since $X$ is compact) is a finite union of compactopen sets of the form $a+d p{ }_{p}{ }_{p}$. (Warning: Not all open sets are compact, for example, $x-\{0\}$.)

Definition. An $\Omega_{p}$-valued measure $\mu$ on $X$ is a finitely additive bounded map from the set of compact open UCX to $\Omega_{p}$.

If we are given the values of a function $\mu$ only on the sets $+d p{ }^{N} Z_{p}$, such a $\mu$ extends to a measure on all compact-open $U$ if and only if these values are bounded and for all a

$$
\begin{equation*}
\mu\left(a+d p{ }_{p}\right)=\sum_{b=0}^{p-1} \mu\left((a+b d p)+d p^{N+1} Z_{p}\right), \tag{2.2}
\end{equation*}
$$

i.e., we need only check additivity for the disjoint unions (2.1).

An equivalent definition of a measure is: a bounded linear functional $\mathrm{f} \longmapsto \int \mathrm{fd} \mu$ on the $\Omega_{\mathrm{p}}$-vector space of locally constant functions on $X$ (i.e., functions which are a finite linear combination of characteristic functions of compact-open sets).

A routine verification shows that, if $f: X \rightarrow \Omega_{p}$ is any continuous function, and we write $f$ as a uniform limit of locally constant functions $f_{i}$, then the limit of the Riemann sums $\int f_{i} d \mu$ exists and depends only on $f: \int f d \mu=\lim \int f_{i} d \mu$. For example, we can evaluate $\int \mathrm{fd} \mu$ as the limit

$$
\begin{equation*}
\int_{\mathrm{fd} \mu}=\lim _{N \rightarrow \infty} \sum_{0 \leq a<d \mathrm{p}} \mathrm{~N} f(a) \mu\left(a+d \mathrm{p}_{\mathrm{p}}^{Z_{\mathrm{P}}}\right) \tag{2.3}
\end{equation*}
$$

Clearly, the $\Omega_{p}$-valued measures form an $\Omega_{p}$-vector space.
For more detailed proofs, see, for example, [53].
Remark. Much more general p-adic measures have been defined: measures on more general types of p-adic spaces $X$ (Mazur, Manin,

Katz), measures which take values in spaces of modular forms (Katz) or spaces of operators (Vishik), unbounded measures (Manin, Vishik).

Basic example. Fix $z \in \Omega_{p}$ so that $\varepsilon=z^{d}$ is not in
$D_{1}\left(I^{-}\right)$. Then $\left|\varepsilon^{P^{N}}-I\right|_{p} \geq 1$ for all $N$. The most important case
is when $\varepsilon=z^{d}$ is a root of one which is not a $p^{\mathbb{N}}$-th root of one for any $N$. Define

$$
\mu_{z}\left(a+d p N_{p}\right)=\frac{z^{a}}{1-\varepsilon^{p^{N}}} .
$$

This gives a measure, since boundedness is ensured by stipulating that $\varepsilon^{\mathrm{P}^{\mathrm{N}}}$ is not close to 1 , and the verification of additivity reduces to summing a geometric progression:

$$
\begin{aligned}
& \sum_{b=0}^{p-1} \mu_{z}\left(a+b d p^{N}+d p^{N+1} z_{p}\right)=\frac{1}{1-\varepsilon^{p^{N+1}}} \sum_{b=0}^{p-1} z^{a+b d p^{N}} \\
= & \frac{z^{a}}{1-\varepsilon^{p^{N+1}}} \sum_{b=0}^{p-1} \varepsilon^{b p^{N}}=\frac{z^{a}}{1-\varepsilon^{p^{N}}}=\mu_{z}\left(a+d p^{N} z_{p}\right)
\end{aligned}
$$

An especially simple case, considered by Osipov [78], occurs when $d=1$ and $z=\varepsilon$ is a ( $p-1$ )-th root of one, in which case the denominator $1-\varepsilon^{p^{N}}=1-\varepsilon$ is simply a constant.

Since the space $X$ "brings together" $Z / d Z$ and $Z_{p}$, we have two natural sources of continuous functions on $X$. (I) Any $\mathrm{f}: \mathrm{Z} \longrightarrow \Omega_{\mathrm{p}}$ having period d can be considered as a continuous (in fact, locally constant) function on $X$ by setting $f(x)=f(a)$ for $x \in a+d Z_{p}$. (2) Any continuous $f: Z_{p} \longrightarrow \Omega_{p}$ can be pulled back to $X$ by means of the map from $X$ to $Z_{p}$ which "forgets mod d information" (i.e., the map which is the inverse liuit of the projections reduction mod $\mathrm{p}^{\mathrm{N}}: \mathrm{Z} / \mathrm{dp}^{\mathrm{N}} \mathrm{Z} \longrightarrow \mathrm{Z} / \mathrm{p} \mathrm{N}_{\mathrm{Z}}$ ).

We shall look at the following example of the second type of continuous function on $X$. Let $t \in \Omega_{p}$ be any small fixed value (namely, ord $t>1 /(p-1)$ ). Then $e^{i x}=\Sigma t^{i} x^{i} / i!$ is a continuous function on $X$; its value at an $X \in X$ is determined by approxi-
mating $x$ by $a$ for which $x \in a+d p^{N} Z_{p}$.
Now let $£$ be a function of period $d$, and consider the function $e^{t x} f(x)$ on $x$. We can integrate this function using (2.3) and summing the geometric progression. We shall write $d \mu_{z}(x)$ to remind ourselves that $x$ (and not $t$ ) is the variable of integra-

$$
\begin{aligned}
& \text { tion. We have: } \\
& \qquad \begin{aligned}
\int e^{t x_{f}(x) d \mu_{z}(x)} & =\lim _{N \rightarrow \infty} \frac{1}{1-\varepsilon^{p^{N}}} \sum_{0 \leq a<d p}^{N} e^{a t} f(a) z^{a} \\
& =\sum_{a=0}^{d-1} f(a) z^{a} e^{a t} \lim _{N \rightarrow \infty} \frac{1}{1-\varepsilon^{p^{N}}} \sum_{b=0}^{p^{N}-1}\left(z e^{t}\right)^{b d} \\
& =\sum_{a=0}^{d-1} \frac{f(a) z^{a} e^{a t}}{1-\varepsilon e^{d t}} \lim _{N \rightarrow \infty} \frac{1-\varepsilon^{p^{N}} e^{d p^{N} t}}{1-\varepsilon^{p^{N}}}
\end{aligned}
\end{aligned}
$$

Since $e^{d p^{N} t}$ approaches 1 as $N \rightarrow \infty$, the limit is 1 , and we obtain

$$
\begin{equation*}
\int e^{t x_{f(x)}^{a i n}} d \mu_{z}(x)=\sum_{a=0}^{d-1} \frac{f(a) z^{a} e^{a t}}{1-\varepsilon e^{d t}} \tag{2.4}
\end{equation*}
$$

Notice that the right side of (2.4) is the same function that appeared in the expression for $L(-k, f, z)$ in 81 , except that in (1.4) the values of $z$ and $f$ were complex, while in (2.4) they are p-adic. The most important case of (1.4) occurs when $f$ takes algebraic values, for example, when $f=\chi:(Z / d Z) * \rightarrow C^{*}$ is a character. Thus, suppose that in (1.4) both $z$ and the values of $f$ are contained in a finite extension $K$ of $Q$. If we imbed $K$ in $\Omega_{p}$, we can identify $z$ and $f(a)$ simultaneousiy as complex or as p-adic numbers.

To construct such an imbedding, choose any prime ideal $P$ of $K$ dividing $p$. Introduce the "P-adic" topology on $K$ in the same way as the p-adic topology was introduced on $Q: x \in K$ is considered to be small if the fractional ideal (x) is divisible by a large positive power of $P$. Then complete $K$ in this topology. Since $P \mid(p)$, the resulting complete field $K_{P}$ contains $Q_{p}$, and
is an algebraic extension of $Q_{p}$. For more details, see [59]. In what follows, we shall suppose that such an imbedding ${ }^{1} p: K \hookrightarrow \Omega_{p}$ has been chosen once and for all, so that any expression involving complex numbers which all lie in $K$ can be simultaneously viewed as a p-adic expression.

In particular, $\mathrm{L}(-\mathrm{k}, \mathrm{f}, \mathrm{z})$ can be considered p-adically. Then, comparing (2.4) with (1.4), we obtain

$$
\begin{aligned}
\sum_{k=0}^{\infty} L(-k, f, z) \frac{t^{k}}{k!} & =\int e^{t x_{f}(x) d \mu_{z}(x)} \\
& =\sum_{k=0}^{\infty} \int x^{k} f(x) d \mu_{z}(x) \frac{t^{k}}{k!} .
\end{aligned}
$$

Since this holds for all $t$ with ord ${ }_{p}>1 /(p-1)$, we can equate coefficients and obtain

$$
\begin{equation*}
\mathrm{L}(-\mathrm{k}, \mathrm{f}, \mathrm{z})=\int \mathrm{x}^{\mathrm{k}} \mathrm{f}(\mathrm{x}) \mathrm{d} \mu_{\mathrm{z}}(\mathrm{x}) . \tag{2.5}
\end{equation*}
$$

As an application of (2.5), one can now study p-adically the values at -k of the Riemann zeta function, since, if we take any positive integer $r$ prime to $p$, we have $*$

$$
\sum_{\varepsilon^{r}=1, \varepsilon \neq 1} L(s, 1, \varepsilon)=\sum_{n=1}^{\infty} n^{-s}\left\{\begin{array}{ll}
r-1 & \text { if } r \mid n \\
-1 & \text { if }\left.r\right|_{n}
\end{array}=\left(r^{1-s}-1\right) \zeta(s)\right.
$$

Thus, for $d=1, X=Z_{p}$, and $\mu$ defined as the sum of $\mu_{\varepsilon}$ over all $\varepsilon$ with $\varepsilon^{r}=\stackrel{p}{1}, \varepsilon \neq 1$, we have
$\zeta(-k)=\frac{1}{r^{k+1}-1} \int x^{k} d \mu(x)$.
Remark. The relation between this $\mu$ and Mazur's measures $\mu_{\alpha}$ (see [53]) is that $\mu=\mu_{\text {Mazur, } \alpha}$ for $\alpha=1 / r$.
3. p-adic interpolation

For simplicity, we first treat the case of the Riemann zeta function, and take $d=1, X=z_{p}$. We know that the values

$$
\begin{equation*}
\zeta(1-k)=-\frac{1}{k} B_{k}=\frac{1}{r^{k}-1} \cdot \int^{\mathrm{p}} \mathrm{x}^{\mathrm{k}-\mathrm{I}} \mathrm{~d} \mu(\mathrm{x}) \tag{3.1}
\end{equation*}
$$

are rational numbers (we have replaced $k$ by $k-1$ in (2.6)). It would be nice to find a continuous p-adic function $\zeta_{p}: z_{p} Q_{p}$
which agrees with $\zeta$ on all $1-k$. Since the set $\{1-k\}$ is dense in $z_{p}$, there can be at most one such $\zeta_{p}$. Such a $\zeta_{p}$ exists if and only if

$$
\begin{aligned}
\mathrm{k}_{1} \text { close to } \mathrm{k}_{2} & \stackrel{p-a d i c a l l y}{ } \\
& \Longrightarrow \zeta\left(1-\mathrm{k}_{1}\right) \text { close to } \zeta\left(1-\mathrm{k}_{2}\right) \text { p-adically }
\end{aligned}
$$

This is not the case, however, and we must first modify the zeta function.

We define a new complex analytic function by setting $\zeta^{*}(\mathrm{~s})=\sum_{\mathrm{p} / \mathrm{n}} \mathrm{n}^{-\mathrm{s}}=\left(1-\mathrm{p}^{-\mathrm{s}}\right) \zeta(\mathrm{s})$
for $R e s>1$ (and for other $s \in C$ by analytic continuation). $\zeta^{*}$ is obtained from $\zeta$ in a similar way to how the Artin-Hasse exponential was obtained from $e^{x}$ in §I. 3 (see the identity (3.2) in Chapter I; the terms with $\mathrm{p} \mid \mathrm{n}$ are omitted to define $\mathrm{E}_{\mathrm{p}}(\mathrm{x})$ ). This procedure is often called "removing the Euler factor at $p$ ", because

$$
\zeta^{*}(s)=\left(1-p^{-s}\right) \zeta(s)=\left(1-\mathrm{p}^{-s}\right) \Pi_{\ell} \frac{1}{1-\ell^{-s}}=\prod_{\ell \neq \mathrm{p}} \frac{1}{1-\ell^{-s}} .
$$

There is yet another way to view $\zeta^{*}$. Let us return to the measures $\mu_{\varepsilon}$. It is easy to see that there does not exist a trans-Iation-invariant (bounded) p-adic measure, i.e., a $\mu$ on $Z_{p}$ such that

$$
\mu\left(a_{1}+p^{N_{Z}}\right)=\mu\left(a_{2}+p_{p}^{N_{p}}\right) \text { for all } a_{1}, a_{2}
$$

However, the measures $\mu_{\varepsilon}$ on $z_{p}$ (for any $\varepsilon \notin D_{1}\left(1^{-}\right)$) have the closest possible property, namely:

$$
\mu_{\varepsilon^{p}}^{\left(a+p^{N} Z_{p}\right)}=\mu_{\varepsilon}\left(a p+p^{\left.N+I_{Z_{p}}\right),}\right.
$$

as follows trivially from the definition. This implies that for any continuous function $f$ on $Z_{p}$

$$
\begin{equation*}
\int_{p Z_{p}} f(x) d \mu_{\varepsilon}(x)=\int_{Z_{p}} f(p x) d \mu_{\varepsilon^{p}}^{p}(x) \tag{3.2}
\end{equation*}
$$

(where for $U \subset X, \int_{U} f$ of course means $\left.\int(f)\right|_{U}$ extended by zero to $X-U$ ) ). Now let $\varepsilon^{r}=1, p \| r$. Since raising to the $p$-th power permutes $r$-th roots of one, we have (where we again let $\mu=$

$$
\begin{align*}
&\left.\sum \mu_{\varepsilon}\right): \\
& \int_{Z_{p}^{*}} x^{k-1} d \mu(x)=\left(\int_{Z_{p}}-\int_{p Z_{p}}\right) x^{k-1} d \mu(x) \\
&=\int_{Z_{p}} x^{k-1} d \mu(x)-\int_{Z_{p}}(p x)^{k-1} d \mu(x) \\
&=\left(1-p^{k-1}\right) \int_{Z_{p}} x^{k-1} d \mu(x) \\
& \text { Dividing by } r^{k}-1, \text { we obtain by }(3.1)  \tag{3.3}\\
& \frac{1}{r^{k}-1} \int_{Z_{p}^{*}} x^{k-1} d \mu(x)=\left(1-p^{k-1}\right) \zeta(1-k)=\zeta^{*}(1-k) .
\end{align*}
$$

Thus, removing the Euler factor is equivalent to restrictigg the domain of integration from $Z_{p}$ to $z_{p}^{*}$.

Now suppose that two values $k_{1}$ and $k_{2}$ are close p-adically, and are also in the same congruence class mod $p-1$, that is, suppose that $k_{1}-k_{2}=(p-1) p_{m} N_{m}, m \in Z$. Then we compare the integrand in (3.3) for $k_{1}$ and $k_{2}$ :

$$
\frac{x^{k_{1}-1}}{x^{k_{2}-1}}=x^{k_{1}-k_{2}}=\left(x^{p-1}\right)^{p_{m}^{N}}
$$

But for $x \in Z_{p}^{*}, x^{p-1} \equiv 1(\bmod p)$ (because $a^{p-1}=1$ for $a \epsilon$ $F_{p}^{*}$ ), and it is easy to see (using the binomial expansion) that

$$
\begin{equation*}
\left(\mathrm{x}^{\mathrm{p}-1}\right)^{\mathrm{p}^{\mathrm{N}} \mathrm{~m}} \equiv 1\left(\bmod _{\mathrm{k}_{\mathrm{p}}-1}^{\mathrm{N}+1}\right) \tag{3.4}
\end{equation*}
$$

Thus, $x^{k_{1}-1}$ and $x^{k_{2}-1}$ are close together p-adically. Hence,
their integrals over the compact set $Z_{p}^{*}$ are also close together; in fact, it is easy to see that

$$
\int_{\mathrm{Z}_{\mathrm{p}}^{*}} x^{k_{1}-1} d \mu(x) \equiv \int_{z_{\mathrm{p}}^{*}} x^{k_{2}-1} d \mu(x) \quad\left(\bmod p^{N+1}\right)
$$

If we further assume that $k_{1} \not \equiv 0(\bmod p-1)$, and if we take $\underset{k_{i}}{r}$ to be a primitive ( $p-1$ )-th root of one modulo $p$ (so that $p / r^{\left.k_{i}-1\right), ~}$ then we have: $1 /\left(r^{k} 1-1\right) \equiv 1 /\left(r^{k} 2-1\right)\left(\bmod p^{N+1}\right)$. Multiplying these two congruences and using (3.3) and (3.1), we obtain the

Kummer congruences. If $k_{1} \equiv k_{2}\left(\bmod (p-1) p^{N}\right)$ and $(p-1) / k_{1}$, then

$$
\left(1-p^{k_{1}-1}\right) \frac{{ }_{k_{1}}}{k_{1}} \equiv\left(1-p^{k_{2}-1}\right) \frac{{ }^{B_{k_{2}}}}{k_{2}}\left(\bmod p^{N+1}\right),
$$

where both sides of the congruence are rational numbers in $Z_{p}$

## (i.e., without $p$ in the denominator).

Thus, the Kummer congruences, which were originally thought to be merely a number theoretic curiosity, are now seen to arise naturally from the simple fact that: if two functions are close together, then their integrals over a compact set are also close together.

We can now define the p-adic zeta function $\zeta_{p}(s)$ by letting $1-\mathrm{k}$ approach s p-adically, but fixing a class modulo $\mathrm{p}-1$, i.e., fixing $k_{0} \in\{0,1, \ldots, p-2\}$ and only choosing $k$ which are con-

$$
\begin{aligned}
& \text { gruent to } \mathrm{k}_{0}(\bmod \mathrm{p}-1) \text {. Thus, we define } \\
& \zeta_{\mathrm{p}, \mathrm{k}_{0}}(\mathrm{~s})=\lim _{1-k \rightarrow \mathrm{~s}, \mathrm{k}=\mathrm{k}_{0}(\bmod \mathrm{p}-1)}^{\zeta^{*}(1-\mathrm{k})} \\
& =\lim _{1-k \rightarrow s, k \equiv k_{0}}(\bmod p-1) \frac{1}{r^{k}-1} \int_{Z_{p}^{*}} x^{k-1} d \mu(x)
\end{aligned}
$$

where $\omega$ is the locally constant function on $Z_{p}^{*}$ which takes a $p$-adic integer to the Teichmuller representative of its first digit, and as before $\langle x\rangle=x / w(x) \equiv I(\bmod p)$. (Thus, $\langle x\rangle^{-s}$ is well defined for p -adic s , see (3.4).) $\zeta_{\mathrm{p}}$ is a p -adic function with p-1 "branches" $\zeta_{\mathrm{p}, \mathrm{k}_{0}}$ for $\mathrm{k}_{0}=0,1, \ldots, \mathrm{p}-2$.

Remark. The classical Mellin transform of a measure $\mu=$ $f(x) d x$ is the function

$$
g(s)=\int_{0}^{\infty} x^{s} f(x) d x
$$

For example, the gamma function is defined as the Mellin transform of $e^{-x_{d x} / x}$ :
$e^{-x} d x / x:$
$\Gamma(s)=\int_{0}^{\infty} e^{-x} x^{s-1} d x$.

Thus, $\zeta_{\mathrm{p}}$ is the p-adic "Mellin-Mazur transform" of the measure $\mu=\sum_{\varepsilon^{r}=1, \varepsilon \neq 1}^{p} \mu^{\mu}=\mu_{\text {Mazur }, 1 / r}$. Strangely, the p-adic $\Gamma$-function, which we shall soon study, is not any type of p-adic Mellin transform, so far as we know.
4. p-adic Dirichlet L-functions

A Dirichlet character $X:(Z / d Z) * \longrightarrow C *$ takes values in a finite (cyclotomic) extension $K$ of $Q$. Recall that we can consider $K$ to be imbedded in $\Omega_{p}$ if we choose a prime ideal $P$ of $K$ dividing $p$ and take the completion of $K$ in the P-adic topology: ${ }^{1} p: K \hookrightarrow \Omega_{p}$. We shall still use the letter $X$ for ${ }^{1} p^{\circ} X$, so $X$ denotes either a complex or p-adic valued character.

A Dirichlet character $X$ is said to be primitive of conductor $d$ if there is no character $X^{\prime}:\left(Z / d^{\prime} Z\right)^{*} \rightarrow C^{*}, d^{\prime}$ a proper divisor of $d$, such that $\chi(n)=\chi^{\prime}(n)$ for all $n$ prime to $d$; equivalently, $X$ is primitive if it is not constant on any subgroup $\left\{x \mid x \equiv 1\left(\bmod d^{\prime}\right)\right\}$ in $(z / d z) *$.

If $X_{1}$ and $X_{2}$ are two primitive Dirichlet characters of conductor $d_{1}$ and $d_{2}$, respectively, then $X_{1} X_{2}$ denotes the primitive Dirichlet character such that $X_{1} X_{2}(n)=X_{1}(n) X_{2}(n)$ whenever $n$ and $d_{1} d_{2}$ have no common factor. This is not the same as the character $n \mapsto X_{1}(n) X_{2}(n)$, which is often imprimitive. For example, if $\chi_{2}=\bar{\chi}_{1}$ is the conjugate character, then $X_{1} X_{2}$ is identically 1 , while $X_{1}(n) X_{2}(n)=0$ if g.c.d. $\left(n, d_{1}\right)>0$. Note that the conductor of $X_{1} \chi_{2}$ divides the least common multiple of $d_{1}, d_{2}$.

If $X$ is a primitive Dirichlet character of conductor $d$ with values in $\Omega_{p}$, we let $\chi_{k}=\chi \omega^{-k}$, where $\omega: n \mapsto \omega(n)$ is the Teichmüller character, which has conductor $p$. Clearly, the conductor of $X_{k}$ is $p d$ if $p \mid d$, and is either $d$ or $d / p$ if $\mathrm{p} \mid \mathrm{d}$.

Let $X$ be, a primitive Dirichlet character of conductor $d$. We
use the relation (2.5) for $\mathrm{f}=\chi$ :
$L(1-k, X, z)=\int_{X} x^{k-1} X(x) d \mu_{z}(x)$.
(We have replaced $k$ by $k-1$ in (2.5).) We want p-adically to interpolate this function of $k$, i.e., to let $k$ approach $1-s \in$ $z_{p}$ and get an $\Omega_{p}$-valued function $L_{p}(s, \chi, z)$. To do this, we must first make two modifications: (1) "remove the Euler factor" by restricting the integral to

$$
X^{*} \text { dēf }_{0<a<d p, p ; a} a+d p Z_{p}
$$

( $X^{*}$ is the inverse image of $Z_{p}^{*}$ under the "forget mod dinformation" map); (2) replace $x$ by $\langle x\rangle=x / \omega(x)$ in $x^{k}$. We thus define

$$
\begin{aligned}
L_{p}(I-k, \chi, z) & d \overline{\bar{e}} f \int_{X^{*}}\langle x\rangle^{k} / x \chi(x) d \mu_{z}(x) \\
& =\int_{X^{*}}\langle x\rangle^{k-1} X_{1}(x) d \mu_{z}(x) \\
& =\int_{X^{*}} x^{k-1} X_{k}(x) d \mu_{z}(x) \\
& =\int_{X} x^{k-1} x_{k}(x) d \mu_{z}(x)-\int_{X}(p x)^{k-1} X_{k}(p x) d \mu_{z^{p}}(x)
\end{aligned}
$$

(see (3.2); the argument is the same for $X$ as for $Z_{p}$ ). Bringing the $p$ outside the second integral and using the above expression for $L(1-k, X, z)$, we conclude that

$$
\begin{equation*}
L_{p}(1-k, \chi, z)=L\left(1-k, \chi_{k}, z\right)-p^{k-1} \chi_{k}(p) L\left(1-k, \chi_{k}, z^{p}\right) \text {. } \tag{4.1}
\end{equation*}
$$

We thus have the following
Proposition. For $X$ a character of conductor $d$, the continuous function from $Z_{p}$ to $\Omega_{p}$

$$
L_{p}(s, \chi, z) \text { deef }_{\bar{e} f} \int_{X^{*}}^{p}\langle x\rangle^{-s} \chi_{1}^{p}(x) d \mu_{z}(x)
$$

$$
\text { interpolates the values } L\left(1-k, x_{k}, z\right)-p^{k-1} \chi_{k}(p) L\left(1-k, x_{k}, z^{p}\right)
$$

This proposition can be used to prove the following theorem.
Theorem (Kubota-Leopoldt [58] and Iwasawa [41]). There exists a unique p-adic continuous (except for a pole at 1 when $X$ is
the trivial character) function $L_{p}(s, X), s \in Z_{p}$, such that
$L_{p}(1-k, X)=\left(1-X_{k}(p) p^{k-1}\right) L\left(1-k, X_{k}\right)$.
Proof. Let $r>1$ be an integer prime to $p d$, and let $z^{r}=$ $1, z \neq 1$. We first note that the ordinary classical l-function can be recovered from the "twisted" L-function $L(s, \chi, z)$ by means of the following relation, which follows immediately from the definitions:

$$
\sum_{z^{r}=1,} \sum_{z \neq 1} L(s, X, z)=\left(r^{1-s} X(r)-1\right) L(s, X), \quad s \in C
$$

Using this for $s=1-k$ and suming (4.1) over nontrivial $r-t h$ roots of unity $z$ (which are only permuted by $z^{H \rightarrow} \rightarrow z^{p}$ ), we have

$$
\begin{aligned}
& \sum_{z^{r}=1, z_{z \neq 1}} L_{p}(1-k, \chi, z)=\left(1-\chi_{k}(p) p^{k-1}\right) \sum_{z^{r}=1, z \neq 1} L\left(1-k, \chi_{k}, z\right) \\
& =\left(r^{k} X_{k}(r)-1\right)\left(1-X_{k}(p) p^{k-1}\right) L\left(1-k, X_{k}\right) \\
& =\left(\langle r\rangle^{k} \chi^{\prime}(r)-1\right)\left(1-X_{k}(p) p^{k-1}\right) L\left(1-k, X_{k}\right) \text { (4.4) }
\end{aligned}
$$

## So we define

$$
\begin{align*}
L_{p}(s, \chi) & =\frac{1}{d e f} \frac{1}{\langle r\rangle^{1-s} \chi(r)-1} \quad \sum_{z^{r}=1,} L_{z \neq 1} L_{p}(s, \chi, z) \\
& =\frac{1}{\langle r\rangle^{1-s} \chi(r)-1} \int_{X^{*}}\langle x\rangle^{-s} \chi_{1}(x) d \mu(x), \tag{4.5}
\end{align*}
$$

where $\mu$ is the sum of $\mu_{z}$ over all $z$ with $z^{r}=1, \quad z \neq 1$. The equality (4.2) in the theorem now follows from (4.4). The continuity of $\mathrm{L}_{\mathrm{p}}(\mathrm{s}, \mathrm{X})$ (more precisely, local analyticity) follows because we are taking the integral of a continuous (actually, analytic) function of $s$ and then dividing by an expression which can only vanish if $s=1$ and $X(r)=1 ; r$ can be chosen so that $\chi(r) \neq 1$ unless $X$ is trivial. This concludes the proof of the theorem.

Notice that the function $\zeta_{p, k}$ ( $s$ ) we defined in $\$ 3$ is precisely $L_{p}\left(s, w^{k} 0\right)$. Also note that, while it was necessary to choose $r$ in order to construct both $\zeta_{p, k_{0}}(s)$ and $L_{p}(s, \chi)$, these functions are in fact independent of $r$.
5. Leopoldt's formula for $L_{p}(1, x)$.

Recall [13] the classical formula for $L(1, \chi)$, which can be derived by Fourier inversion on the group $G=Z / d Z$. Let $\zeta$ be a fixed primitive $d$-th root of 1 , and define.

$$
\hat{\mathrm{f}}(\mathrm{a})=\sum_{\mathrm{b} \in \mathrm{G}} \mathrm{f}(\mathrm{~b}) \zeta^{-\mathrm{ab}}
$$

for a function $f$ on $G$. Then

$$
\begin{equation*}
f(b)=\frac{1}{\mathrm{~d}} \sum_{a \in G} \hat{f}(a) \zeta^{a b} . \tag{5.1}
\end{equation*}
$$

Applying Fourier inversion (5.1) to $f_{s}(b)=\Sigma_{n=b(\bmod d)} n^{-s}$ (suppose Res $\mathrm{s}>1$ ) and using the definition of $\mathrm{L}(\mathrm{s}, \mathrm{X})$ and $\mathrm{L}(\mathrm{s}, 1, \mathrm{z})=\Sigma \mathrm{z}^{\mathrm{n}} \mathrm{n}^{-s}$, we have

$$
\begin{aligned}
L(s, \chi) & =\sum_{0 \leq b<d} \chi(b) f_{s}(b)=\frac{1}{d} \sum_{a, b} \chi(b) \hat{f}_{s}(a) \zeta^{a b} \\
& =\frac{1}{\bar{d}} \sum_{j} \chi(j) \zeta^{j} \Sigma_{a} \bar{\chi}(a) \hat{f}_{s}(a) \quad \text { (where } j=a b \text { ) } \\
& =\frac{g_{X}}{d} \sum_{a} \bar{\chi}(a) L\left(s, 1, \zeta^{-a}\right)
\end{aligned}
$$

where $g_{\chi}=\Sigma \chi(j) \zeta^{j}$ is the Gauss sum. Letting $s \rightarrow 1$ and noting that $L(1,1, z)=-\log (1-z)$, we obtain

$$
\begin{equation*}
L(1, x)=-\frac{g_{X}}{d} \sum_{0<a<d} \bar{\chi}(a) \log \left(1-\zeta^{-a}\right) \tag{5.2}
\end{equation*}
$$

We now proceed to the p-adic case.
Theorem (Leopoldt [64]).

$$
\begin{equation*}
L_{p}(1, \chi)=-\left(1-\frac{X(p)}{p}\right) \frac{g_{X}}{d} \sum_{0<a<d} \bar{\chi}(a) \ln _{p}\left(1-\zeta^{-a}\right) . \tag{5.3}
\end{equation*}
$$

The purpose of this section is to prove this theorem.
Note that (5.3) differs from (5.2) in two respects: the expected "removal of the Euler factor", giving the term ( $1-\chi(\mathrm{p}) / \mathrm{p}$ ); and the replacement of $\log$ by $1 n_{p}$. The validity of the padic formula (5.3) might seem surprising at first because of the replacement of $\log \left(1-\zeta^{-a}\right)$ by $1 n_{p}\left(1-\zeta^{-a}\right)$, since the formal series for the former does not converge p-adically, i.e., we need properties (2) and (3) of the proposition in $§ 1.3$ in order to evaluate
$\ln _{p}\left(1-\zeta^{-a}\right)$. But the proof of Lemma 1 below shows that the same formal series, along with the p-adic version of analytic continuation, really do lie behind Leopoldt's formula, despite first impressions.

Lemma 1. If $z \& D_{1}\left(1^{-}\right)$and $\mu_{z}$ is the measure on $Z_{p}$ given by $\mu_{z}\left(a+p^{N} Z_{p}\right)=z^{a} /\left(1-z^{N}\right)$, then

$$
\int_{z_{p}^{*}}^{\infty} \frac{1}{x} d \tilde{z}_{z}(x)=-\frac{1}{p} \ln _{p} \frac{(1-z)^{p}}{1-z^{P}}
$$

Proof. If $|z|_{\mathrm{p}}<1$, then the left side is (the ' denotes omission of indices divisible by $p$ ):

$$
\begin{aligned}
\lim _{N \rightarrow \infty} \sum_{0<j<p^{N}}^{\prime} \frac{z^{j}}{j} \frac{1}{1-z^{p^{N}}} & =\lim _{N \rightarrow \infty}\left(\sum_{0<j<p} \frac{z^{j}}{j}-\sum_{0<j<p} N-1 \frac{z^{p j}}{p j}\right) \\
& =-\frac{1}{p}\left(\ln _{p}(1-z)^{p}-\ln _{p}\left(1-z^{p}\right)\right) .
\end{aligned}
$$

We now use analytic continuation to extend the equality from $|z|_{p}<1$ to all $z \in D_{1}\left(I^{-}\right)$. As we remarked at the end of ${ }^{5} .3$, a function is said to be Krasner analytic on the complement of $D_{1}\left(I^{-}\right)$if it is a uniform limit of rational functions with poles in $D_{1}\left(1^{-}\right)$. The basic fact we need about such functions (see [57]) is that if two Krasner analytic functions on the complement of $D_{1}\left(1^{-}\right)$are equal on a disc, then they are equal everywhere on the complement of $D_{1}\left(1^{-}\right)$. Thus, if we show that the two sides of Lemma 1 are each Krasner analytic functions of $z$ on the complement of $D_{1}\left(1^{-}\right)$, then, since they are equal on the disc $|z|_{p}<1$, they must be equal for all $z \notin D_{1}\left(1{ }^{\prime \prime}\right)$.

Note that by writing

$$
\begin{aligned}
& \quad \frac{(1-z)^{p}}{1-z^{p}}= \begin{cases}1+\frac{1}{1-z^{p}} \sum_{0<j<p}\binom{p}{j}(-z)^{j}, & \text { if } \quad|z|_{p} \leq 1, \quad z \notin D_{1}\left(1^{-}\right) ; \\
1+\frac{1}{I-z^{-p}} \sum_{0<j<p}\binom{p}{j}(-z)^{-j}, & \text { if }|z|_{p}>1,\end{cases} \\
& \text { we see that } z \notin D_{1}\left(1^{-}\right) \Longrightarrow(1-z)^{p} /\left(1-z^{p}\right) \in D_{1}\left(I^{-}\right) \quad \text { (in fact, }
\end{aligned}
$$

its distance from 1 is $\leq 1 / \mathrm{P}$ ). Hence, the right side of the equality in the lemma is the uniform limit of the rational functions (with poles in $\mathrm{D}_{1}\left(1^{-}\right)$)
$\frac{1}{p} \sum_{j=1}^{N} \frac{(-1)^{j}}{j}\left(\frac{(1-z)^{p}}{1-z^{p}}-1\right)^{j}$,
and the left side is the uniform limit of the rational functions (with poles in $\mathrm{D}_{1}\left(\mathrm{l}^{-}\right)$)

$$
\sum_{0<j<p}^{N} \frac{z^{j}}{j} \frac{1}{1-z^{p^{N}}}
$$

This concludes the proof of the lemana.
Lemma 1 will be applied when $z$ is a root (but not a $p^{N}$-th root) of unity.

Remarks. 1. If $z$ is a ( $p-1$ )-th root of 1 , the right side of Lemma 1 becomes $-(1-1 / p) \ln _{p}(1-z)$. For example, setting $z=-1$ gives the following p-adic limit for $\ln _{\mathrm{p}} 2$ :

$$
\ln _{p} 2=-\frac{p}{2(p-1)} \lim _{N \rightarrow \infty} \sum_{0<j<p} \frac{(-1)^{j}}{j}
$$

2. Lemma 1 is the key step in our proof of Leopoldt's formula for $L_{p}(1, \chi)$. As mentioned before, the subtlety in Leopoldt's formula is that $\ln _{p}(1-z)$ is not given by the same formal series as $\log (1-z)$, since $z$ is outside the disc of convergence of $I_{p}(1-z)$. However, Lemma 1 shows that if we "correct by the Frobenius" in the Dwork style (see, e.g., [28]), i.e., if we replace (1-z) by ( $1-z)^{\mathrm{P}} /\left(1-z^{\mathrm{p}}\right)$, then the resulting series is globally analytic out to roots of unity. (We shall see further examples of "correcting by the Frobenius", e.g., in §6.) The effect of this step on the formula for $L_{p}(I, X)$ is to bring out the Euler factor ( $1-\mathrm{X}(\mathrm{p}) / \mathrm{p}$ ), as we shall see below (in (5.4)).

The other ingredient in the proof of Leopoldt's formula is the analog of the Fourier inversion (5.i) used in the classical case.

Lemma 2. Suppose that $\chi$ is a primitive Dirichlet character $\bmod d, \quad \zeta$ is a fixed primitive d-th root of unity, $g_{x}=\Sigma \chi(j) \zeta^{j}$, $z \neq 1$ is an r-th root of unity, where $r$ is prime to $p d$, and $\mathrm{f}: \mathrm{x} \rightarrow \Omega_{\mathrm{p}}$ is any continuous function. Then

$$
\int_{X} X f d \mu_{z}=\frac{g_{X}}{d} \sum_{0 \leq a<d} \bar{\chi}(a) \int_{X} f d \mu_{\zeta^{-a} z} .
$$

To prove Leman 2, by linearity and continuity of both sides it suffices to prove it when $f$ is the characteristic function of
$j+d p^{N} Z_{p}$, i.e., to prove that $\chi(j) z^{j} /\left(1-z^{d p^{N}}\right)=$
$\frac{g_{\chi}}{d} \sum \bar{\chi}(a) z^{j} \zeta^{-a j} /\left(1-z^{d p}\right)$. But this reduces to: $g_{\chi} \bar{g}_{\chi}=d$.
Proof of the theorem. We first prove an analogous formula for the "twisted" $\mathrm{L}_{\mathrm{p}}(1, \chi, z)$.

Note that, if $f: X \rightarrow \Omega_{p}$ comes from pulling back a function (also denoted $f$ ) on $z_{p}$ using the projection ("forget mod d information") from $x$ to $Z_{p}$, then we can replace $X$ by $Z_{p}$ in $\int_{X} f d \mu_{z}$, where $\mu_{z}$ on $Z_{p}$ is defined by the same formula as on $X$ with $d$ replaced by $1: \mu_{z}\left(j+p^{N} Z_{p}\right)=z^{j} /\left(1-z^{p^{N}}\right)$. To see this, one reduces to the case when $f$ is the puili-back of the characteristic function of $j+p^{N} Z_{p}$, which is checked easily.

Applying Lemma 2 and the preceding remark to the function $f(x)=\frac{1}{x} \cdot\left(\right.$ characteristic function of $\left.X^{*}\right)$, we obtain

$$
\begin{align*}
L_{p}(1, \chi, z) & =\int_{X^{*}} \frac{X(x)}{x} d \mu_{z}=\frac{g_{X}}{d} \sum_{0 \leq a<d} \bar{X}(a) \int_{Z_{p}^{*}} \frac{1}{x} d \mu_{\zeta^{-a} z} \\
& =-\frac{g_{X}}{d} \sum_{0 \leq a<d} \bar{\chi}(a) \frac{1}{p} \ln _{p_{p}} \frac{\left(1-\zeta^{-a} z\right)^{p}}{} \text { by Lemma } 1 \\
& =-\frac{g_{X}}{d}\left(1-\frac{X(p)}{p}\right) \sum_{0<a<d} \bar{x}(a) 1 n_{p}\left(1-\zeta^{-a} z\right) . \tag{5.4}
\end{align*}
$$

Since $r>1$ is any integer prime to $p d$, we may choose $r$ so that $X(r) \neq 1$ and then use (4.5) with $s=1$ to express $L_{p}(1, \chi)$ in terms of the $L_{p}(1, \chi, z)$. We obtain

$$
\begin{aligned}
L_{p}(1, \chi) & =\frac{1}{\chi(r)-1} \sum_{z^{r}=1, z \neq 1} L_{p}(1, \chi, z) \\
& =-\frac{g_{\chi}}{d}\left(1-\frac{\chi(p)}{p}\right)\left[\frac{1}{\chi(r)-1} \sum_{0<a<d} \bar{\chi}(a) \sum_{z=1,} \sum_{z \neq 1} \ln _{p}\left(1-\zeta^{2} z_{z}\right)\right]
\end{aligned}
$$

Since the inner summation is equal to $\ln _{p}\left(1-\zeta^{-a r}\right)-\ln _{p}\left(1-\zeta^{-a}\right)$, the term in square brackets is immediately seen to equal $\Sigma_{0<a<d} \bar{X}(a) \ln _{p}\left(1-\zeta^{-a}\right)$, as desired.

This completes the proof of Leopoldt's formula
We shall later prove one more formula for the p-adic l-function, which relates its behavior at 0 to the p-adic ganina function But first we take up the p-adic gamma and log gamma functions.
6. The p-adic gamma function

First recall some properties of the classical gamna function:
(1) It is a meromorphic function on the complex plane with poles at $0,-1,-2,-3, \ldots$.
(2) $\Gamma(x+1)=x \Gamma(x), \quad \Gamma(k)=(k-1)!$
(3) $\Gamma(x) \cdot \Gamma(1-x)=\frac{\pi}{\sin (\pi x)}=\frac{2 \pi i}{e^{2 \pi i x}-1} e^{\pi i x}$
(4) $\Gamma(x)=\lim _{n \rightarrow-\infty} \frac{1 \cdot 2 \cdots(n-1)}{x(x+1) \cdots(x+n-1)} n^{x}$
(5) Gauss multiplication formula: for $\mathrm{m}=1,2, \ldots$
$\prod_{h=0}^{m-1} \Gamma\left(\frac{x+h}{m}\right)=(2 \pi)^{(m-1) / 2} m^{\frac{1}{2}-x} \Gamma(x)$,
for example ( $\mathrm{x}=\mathrm{I}$ )

$$
\prod_{h=1}^{m-1} \Gamma\left(\frac{h}{m}\right)=(2 \pi)^{(m-1) / 2 m^{-1 / 2},}
$$

o that if we divide these two equations we obtain

$$
\begin{equation*}
\frac{\prod_{h=0}^{m-1} \Gamma\left(\frac{x+h}{m}\right)}{\Gamma(x) \prod_{h=1}^{m-1} \Gamma\left(\frac{h}{m}\right)}=m^{1-x} . \tag{6.1}
\end{equation*}
$$

We now proceed to the p-adic theory. For simplicity, we shall assume that $p \neq 2$. (Minor modifications are sometimes needed when $\mathrm{p}=2$.)

Proceeding naively, we would like to construct a function $\Gamma_{p}(s)$ on $Z_{p}$ which interpolates $\Gamma(k)=\Pi_{j<k} j$, i.e., so that $\Gamma(k)$ approaches $\Gamma_{p}(s)$ as $k$ runs through any sequence of positive integers which approaches $s$ p-adically. However, $\Gamma(k)$ is divisible by a large power of $p$ for large $k$; hence $\Gamma_{p}(s)$ would have to be identically zero, which is useless. So, as in the case of the zeta function, we must modify the values $\Gamma(k)$

We might improve the situation if we eliminate the $j$ 's which are divisible by $p$, much as we "removed the Euler factor at $p$ " from $\zeta(s)$ to get $\zeta^{*}(\mathrm{~s})$. Thus, suppose we take $\Gamma^{*}(\mathrm{k})=$
$\Pi_{j<k, p / j} j$, which is an integer prime to $p$. Now we can find a continuous function $\Gamma_{p}$ on $Z_{p}$ which agrees with $\Gamma *$ on positive integers if and only if
$\mathrm{k}_{1}$ close to $\mathrm{k}_{2} \mathrm{p}$-adically $\Longrightarrow \Gamma *\left(\mathrm{k}_{1}\right)$ close to $\Gamma *\left(\mathrm{k}_{2}\right)$ i.e., if and only if. $\Gamma^{*}\left(k_{2}\right) / \Gamma^{*}\left(k_{1}\right) \equiv 1\left(\bmod p^{N}\right)$ with $N$ large whenever $k_{2}-k_{1}$ is highly divisible by $p$. To check this, take for example $k_{2}=k_{1}+p^{n}$. But it is easy to show that in that case the quotient $\Gamma^{*}\left(\mathrm{k}_{2}\right) / \Gamma^{*}\left(\mathrm{k}_{1}\right)=\Pi_{\mathrm{k}_{1} \leq \mathrm{j}<\mathrm{k}_{2}}, \mathrm{p} / \mathrm{j} \mathrm{j}$ is congruent to $-1 \bmod p^{n}$. (Namely, in any finite abelian group $G$ we have $\Pi_{g \in G} g=\Pi_{g=g^{-1}} g$; apply this to $G=\left(Z / p^{n} Z\right) *$.) So the sign is wrong, and we have to make one final modification. We define

$$
\begin{equation*}
\Gamma_{p}(k)=(-1)^{k} \prod_{j<k, p / j} j, \quad \Gamma_{p}(s)=\lim _{k \rightarrow s} \Gamma_{p}(k) . \tag{6.2}
\end{equation*}
$$

Using the generalized Wilson's theorem cited above:

$$
\prod_{k \leq j<k+p^{n}, p \nmid j} j \equiv-1 \quad\left(\bmod p^{n}\right),
$$

it is simple to check that the limit in (6.2) exists, is independent of how $k$ approaches $s$, and determines a continuous function
on $Z_{p}$ with values in $Z_{p}^{*}$.
We verify some basic properties of $\Gamma_{p}$ :
(1) $\Gamma_{p}(0)=1$, and $\frac{\Gamma_{p}(x+1)}{\Gamma_{p}(x)}=\left\{\begin{array}{ccc}-x & \text { if } & x \in Z_{p}^{*} ; \\ -1 & \text { if } & x \in p Z_{p} .\end{array}\right.$

To prove the second equality, since both sides are continuous on $Z_{p}$, it suffices to prove equality for the positive integers $x=k$ (which are dense in $Z_{p}$ ), and then it is obvious from (6.2). Using (6.2) and this equality, we can compute the first few values $\Gamma_{p}(2)=1, \quad \Gamma_{p}(1)=-1, \quad \Gamma_{p}(0)=1$.
(2) For $x \in Z_{p}$, write $x=x_{0}+p x_{1}$, where $x_{0} \in\{1, \ldots, p\}$
is the first digit in $x$, unless $x \in p Z_{p}$, in which case $x_{0}=p$. Then we have

$$
\begin{aligned}
& \\
& \Gamma_{p}(x) \cdot \Gamma_{p}(1-x)(-1)^{x_{0}} .
\end{aligned}
$$

In fact, to show that the continuous function $f(x)=(-1){ }^{x_{0_{p}}}(x)$. $\Gamma_{p}(1-x)$ equals 1 on $Z_{p}$, it suffices to show that $f(k)=1$ for positive integers $k$. Clearly $f(1)=1$, and a simple verification using property ( 1 ) shows that $f(k+1) / f(k)=1$.
(3) For any positive integer $m, p \nmid m$, we have (here $x_{0}$ and $x_{1}$ are as in property (2)):

$$
\frac{\prod_{h=0}^{m-1} \Gamma_{p}\left(\frac{x+h}{m}\right)}{\Gamma_{p}(x) \prod_{h=1}^{m-1} \Gamma_{p}\left(\frac{h}{m}\right)}=m^{1-x_{0}}\left(m^{-(p-1)}\right)^{x_{1}}
$$

(Note: Since $m^{-(p-1)} \equiv 1(\bmod p)$, it follows that $\left(m^{-(p-1)}\right)^{x_{1}}$ is a well-defined function of p-adic $x_{1}$.) To prove property (3), let $f(x)$ be the left side and $g(x)$ the right side of (6.3). $f$ and $g$ are continuous, and $f(1)=1=g(1)$. Next,

$$
\frac{f(x+1)}{f(x)}=\frac{\Gamma_{p}(x)}{\Gamma_{p}(x+1)} \cdot \frac{\Gamma_{p}\left(\frac{x}{m}+1\right)}{\Gamma_{p}(x / m)}= \begin{cases}1 / m & \text { if } x \in Z_{p}^{*} ; \\ 1 & \text { if } x \in p Z_{p},\end{cases}
$$

while

$$
\frac{g(x+1)}{g(x)}=\left\{\begin{array}{l}
1 / \dot{m} \text { if } x \in Z_{p}^{*} ; \text { since then }(x+1)_{0}=x_{0}+1,(x+1)_{1}=x_{1} \\
1
\end{array} \quad \text { if } x \in p Z_{p}, \text { since then } \begin{array}{l}
(x+1)_{0}=x_{0}-(p-1), \\
(x+1)_{1}=x_{1}+1 .
\end{array}\right.
$$

This proves property (3).
We now discuss an interesting special case of property (3).
Suppose that $\quad x=\frac{r}{p-1} \quad$ is a rational number between 0 and 1 whose denominator divides $p-1$. Then $x_{0}=p-r, x_{i}=\frac{r}{p-1}-1=$ $\left(1-x_{0}\right) /(p-1)$. Note that the left side $f(x)$ of (6.3) is congruent to $\mathrm{m}^{1-\mathrm{x}_{0}}=\mathrm{m}^{(1-\mathrm{x})(1-\mathrm{p})}$ mod p . In addition,

$$
\mathrm{f}(\mathrm{x})^{\mathrm{p}-1}=\left(\mathrm{m}^{\mathrm{p}-1}\right)^{\left(1-\mathrm{x}_{0}-\mathrm{x}_{1}(\mathrm{p}-1)\right)}=1 .
$$

Thus, $f(x)$ is the ( $p-1$ )-th root of 1 congruent to $m^{(1-x)(1-p)}$ mod p ;
$f(x)=\omega\left(m^{(1-x)(1-p)}\right)$.
Now the classical expression $f_{c 1}(x)$ which is obtained from $f(x)$ by replacing $\Gamma_{p}$ by $\Gamma$, is equal to $m^{I-x} \in Q\left(\frac{p-1}{m}\right)$. Let $K^{\prime}=Q(\xi)$, where $\xi$ is a fixed primitive ( $p-1$ )-th root of unity. Then $\operatorname{Gal}\left(\mathrm{K}\left(\frac{\mathrm{p}-1}{-1} \sqrt{\mathrm{n}}\right) / \mathrm{K}\right) \cong \mathrm{Z} /(\mathrm{p}-1) \mathrm{Z}$, with $\sigma_{a}: \frac{\mathrm{p}-1}{} \sqrt{\mathrm{~m}} \rightarrow \xi^{\mathrm{a}} \mathrm{p}-1 \sqrt{\mathrm{~m}}$ for $a \in z /(p-1) z$. Choose a prime ideal $P$ of $K\left(\frac{p-1}{\sqrt{m}}\right)$ which divides $p$. Then $P$ determines an imbedding $l_{p}: K\left(P^{p-1} \sqrt{\mathbb{1}}\right) \hookrightarrow$ $Q_{p}\left(\frac{P^{-1}}{/ n}\right)$. There exists a unique "Frobenius element" Frob $\epsilon$ $\operatorname{Gal}\left(\mathrm{K}\left(\mathrm{P}^{-1} / \sqrt{\mathrm{m}}\right) / \mathrm{K}\right)$ such that $\operatorname{Frob}(\mathrm{x}) \equiv \mathrm{x}^{\mathrm{P}}(\bmod \mathrm{P})$ for every algebraic integer $x$ in $K\left(\frac{p-1}{\sqrt{m}}\right)$.

We then have for $x=r /(p-1)$

$$
f_{c 1}(x)^{1-\text { Frob }}=\frac{m^{1-x}}{\text { Frob } m^{1-x}} \equiv m^{(1-x)(1-p)} \quad(\bmod P) ;
$$

and, since elements of $\operatorname{Gal}\left(\mathrm{K}\left(\mathrm{P}^{\mathrm{P}-1} \sqrt{\mathrm{~m}}\right) / \mathrm{K}\right)$ multiply $\mathrm{m}^{1-\mathrm{X}}$ by roots of unity, we conclude that

$$
i_{p}\left(f_{c i}(x)^{1-F r o b}\right)=\omega\left(m^{(1-x)(1-p)}\right)=f_{p-a d i c}(x)
$$

This phenomenon -- that a classical expression raised to the 1-Frob, where Frob is a p-th power type map, can be identified with a p-adic analog of the classical expression -- occurs in other contexts. For example, let

E: $\quad y^{2}=x(x-1)(x-\lambda), \quad \lambda \in Z$,
be an elliptic curve whose reduction

$$
\bar{E}: \quad y^{2}=x(x-1)(x-\bar{\lambda}), \quad \bar{\lambda} \in z / p z
$$

is nonsingular. Further suppose that $\overline{\mathrm{E}}$ is not "supersingular" (which will be the case if $\bar{\lambda}$ is not a root of the polynomial $\left.\sum_{n=0}^{(p-1) / 2}\binom{(p-1) / 2}{n}^{2} \bar{\lambda}^{n}\right)$. It is known [67,43] that the period $f$ of the holomorphic differential $\mathrm{dx} / \mathrm{y}$ on the Riemann surface (torus) $E$, as a function of the parameter $\lambda$, satisfies the differential equation

$$
\lambda(1-\lambda) f^{\prime \prime}+(1-2 \lambda) f^{\prime}-\frac{1}{4} f^{\prime}=0
$$

whose solution bounded at zero is the hypergeometric series

$$
f(\lambda)=\sum_{n=0}^{\infty}\binom{-1 / 2}{n}^{2} \lambda^{n} \in Q[[\lambda]] .
$$

Although when we consider $f(\lambda)$ as a $p$-adic series it converges only on $D\left(1^{-}\right)$, it turns out that the power series $\theta=f^{1-F r o b}$ defined by $\theta(\lambda)=f(\lambda) / f\left(\lambda^{p}\right)$ converges on $D(\gamma)$ for some $y>1$.

$$
\left.\begin{array}{l}
\text { Now it is well known that the zeta-function of } \overline{\mathrm{E}} \\
\mathrm{Z}\left(\overline{\mathrm{E}} / \mathrm{F}_{\mathrm{p}}\right)=\exp \sum\left(\frac{(\text { number of }}{} \mathrm{F}_{\mathrm{n}^{\text {n }}} \mathrm{points} \text { on } \overline{\mathrm{E}}\right) \\
\mathrm{n}
\end{array}\right) \mathrm{T}^{\mathrm{n}} .
$$

is of the form

$$
Z\left(\stackrel{\rightharpoonup}{E} / F_{p}\right)=\frac{(1-\alpha T)(1-p T / \alpha)}{(1-\widetilde{T})(1-p T)}
$$

where $\operatorname{ord}_{p} \alpha=0$. Dwork [43] proved the following formula for $\alpha$ : $\alpha=\stackrel{\mathrm{P}}{\ominus}(\omega(\bar{\lambda}))$.
Thus, $\alpha$ can be thought of as a sort of "p-adic period".
In Chapter III we shall study another analogy between classical formulas for periods of a curve considered over $C$ and p-adic
formulas for the roots of the zeta function of the curve considered over a finite field.

This concludes our discussion of elementary (easily proved) properties of $\Gamma_{p}$. In Chapter III we shall prove an algebraicity result for certain $\Gamma_{p}$ values and products of values, for example, the algebraicity of $\Gamma_{p}\left(\frac{r}{d}\right)$ if $d \mid p-1$. (More precisely: $\left(\Gamma_{p}\left(\frac{r}{d}\right)\right)^{d}$ $\in Q(\sqrt{d} \sqrt{1})$.) But the proof of this fact uses p-adic cohomology. It would be interesting to find an elementary proof of this algebraicity. After all, the assertion can be stated very simply (here we write $\left.\frac{\mathrm{r}}{\mathrm{d}}=1-\frac{\mathrm{s}}{\mathrm{p}-1}\right)$ :

$$
\lim _{\mathrm{n} \rightarrow \infty} \frac{\left(\mathrm{~s}+\mathrm{sp}+\ldots+s p^{n}+s p^{\mathrm{n}+1}\right)!}{\left(s+s p+\ldots+s p^{n}\right)!p^{s+s p+\ldots+s p^{\mathrm{n}}}} \in \mathrm{Z}_{\mathrm{p}}^{*} \text { is algebraic over } Q \text {. }
$$

For example, the theorem in SIII. 6 will give us the following 5-adic formula:

$$
\Gamma_{5}\left(\frac{1}{4}\right)^{4}=\lim _{\mathrm{n} \rightarrow \infty}\left(\frac{\left(3+3 \cdot 5+\ldots+3 \cdot 5^{\mathrm{n}}+3 \cdot 5^{\mathrm{n}+1}\right)!}{\left(3+3 \cdot 5+\ldots+3 \cdot 5^{n}\right)!5^{3+3 \cdot 5+\ldots+3 \cdot 5^{n}}}\right)^{4}
$$

$$
=3+4 \omega(2)=3+4 \sqrt{-1} \in Z_{5}^{*}
$$

(where $\sqrt{-1}=\omega(2)=2+1 \cdot 5+2 \cdot 5^{2}+1 \cdot 5^{3}+3 \cdot 5^{4}+\ldots \epsilon Z_{5}$ ) and the $\varepsilon_{n} 11$ owing 7 -adic formula:

$$
\begin{aligned}
-\Gamma_{7}\left(\frac{1}{3}\right)^{3} & =\lim _{n \rightarrow \infty}\left(\frac{\left(4+\ldots+4 \cdot 7^{n}+4 \cdot 7^{n+1}\right)!}{\left(4+\ldots+4 \cdot 7^{n}\right)!7^{4+\ldots+4 \cdot 7^{n}}}\right)^{3} \\
& =\frac{-1+3 \sqrt{-3}}{2}=1+3 \omega(4) \in Z_{7}^{*} .
\end{aligned}
$$

No elementary proof is known for either of these equalities.

## 7. The p-adic log gamma function

We start by describing another approach to the p-adic zeta and L-functions, which was the original point of view of Kubota and Leopoldt [58].

It is not hard to prove the following p-adic formula for the
k-th Bernoulli number (see [41]):

$$
\begin{equation*}
B_{k}=\lim _{n \rightarrow \infty} p^{-n} \sum_{0 \leq j<p} n^{k} \tag{7.1}
\end{equation*}
$$

More generally, if $f: Z \rightarrow \Omega_{p}$ has period $d$, and if $B_{k, f}$ is defined for $p$-adic valued $f$ in the same way as for complex valued f (see (1.1)), then we have

$$
\begin{equation*}
B_{k, f}=\lim _{n \rightarrow \infty} \frac{1}{d p^{n}} \sum_{0 \leq j<d p^{n}} f(j) j^{k} . \tag{7.2}
\end{equation*}
$$

The simplest examples of (7.1) are:

$$
\begin{aligned}
& \lim p^{-n} \sum_{j<p} j=\lim p^{-n} \frac{p^{n}\left(p^{n}-1\right)}{2}=-\frac{1}{2}=B_{1} ; \\
& \lim p^{-n} \sum_{j<n}^{n} j^{2}=\lim p^{-n} \frac{p^{n}\left(p^{n}-1\right)\left(2 p^{n}-1\right)}{6}=\frac{1}{6}=B_{2} .
\end{aligned}
$$

This type of Limit $\lim _{n \rightarrow \infty} p^{-n} \sum_{0 \leq j<p} n(j) \quad$ can be used for other $f(x)$ besides $f(x)=x^{k}$.

Definition. Suppose that a subset $\mathrm{V} \subset \Omega_{\mathrm{p}}$ has no isolated points. A function $f: U \rightarrow \Omega$ is called locally analytic if for every $a \in U$ there exist $r \stackrel{p}{a}$ and $a_{i}$ such that for all $x$ in $D_{a}(r) \cap U$

$$
f(x)=\sum_{i=0}^{\infty} a_{i}(x-a)^{i}
$$

It is easy to check that a locally analytic function $£$ can be differentiated in the usual way:

$$
f^{\prime}(x)=\lim _{\varepsilon \rightarrow 0} \frac{f(x+\varepsilon)-f(x)}{\varepsilon}=\sum i a_{i}(x-a)^{i-1} \quad \text { for } x \in D_{a}(r) \cap U
$$

Lemma. If $f$ is locally analytic on $Z_{p}$, then the Iimit $\lim _{n \rightarrow \infty} p^{-n} \sum_{0 \leq j<p}^{n} f(j)$ exists.

To prove the lemna, one easily reduces to the case when $f(x)=$ $\Sigma a_{i} x^{i}$ on $D(1) .\left(T h u s, \quad a_{i} \longrightarrow 0\right.$.) Then we need to show that $\Sigma a_{i} B_{i}$ converges, but this follows because $\left|B_{k}\right|_{p} \leq p$. (More precisely, we have already seen that $B_{k} \in Z_{p}$ if $p-1 / k$, and one can
similarly use the p-adic integral formula

$$
\left(1-p^{k-1}\right) B_{k}=\frac{\mathrm{k}}{1-r^{k}} \int_{Z_{p}} x^{k-1} d \mu(x)
$$

to show that $p B_{k} \in Z_{p}$ if $p-1 \mid k$.)
The approach of Kubota-Leopoldt is based on this lemma and the formulas (7.1) and (7.2). In order to obtain a p-adic function out of (7.1) as $k$ approaches some $p$-adic $s$, we must omit the $j$ which are divisible by $p$ and also restrict to $k \equiv k_{0}(\bmod p-1)$ for some fixed $k_{0}$. If we write $j=\langle j\rangle \omega(j)$ for $p / j$, then we obtain for $k \equiv k_{0}(\bmod p-1) \quad$ (see $\S 3$ for the definition of $\zeta^{*}$ ):

$$
\begin{aligned}
\zeta^{*}(1-k) & =\left(1-p^{k-1}\right)\left(-\frac{B_{k}}{k}\right) \\
& =-\frac{1}{k} \lim _{n \rightarrow \infty} p^{-n}\left(\sum_{0 \leq j<p} n^{k}-p^{k} \sum_{0 \leq j<p^{n-1}} j^{k}\right) \\
& =-\frac{1}{k} \lim _{n \rightarrow \infty} p^{-n} \sum_{\left.0 \leq j<p^{n}, p j_{j}<j\right\rangle^{k} \omega(j)^{k} 0} .
\end{aligned}
$$

We can now define $\zeta_{p, k_{0}}(s)$ by replacing $k$ by $1-s$ and applying the lemma to $f(x)=\left\langle_{x}\right\rangle^{1-s} \omega(x)^{k_{0}}$ (we take $f(x)=0$ on $p Z_{p}$ ). Similarly, the p-adic L-function for a Dirichlet character $\chi$ : $(\mathrm{Z} / \mathrm{dZ}) * \rightarrow \Omega_{\mathrm{p}}^{*} \quad$ can be defined by setting

$$
\begin{equation*}
L_{p}(1-s, x)=-\frac{1}{s}{\underset{n \rightarrow \infty}{ } \lim _{n \rightarrow \infty} \frac{1}{d p^{n}} \sum_{0 \leq j<d p} n, p \|_{j}\langle j\rangle^{s} \chi(j) . ~ . ~ . ~} \tag{7.3}
\end{equation*}
$$

This approach to the construction of $\zeta_{\mathrm{p}, \mathrm{k}_{0}}$ and $\mathrm{L}_{\mathrm{p}}$ can be generalized as follows. Let $X=\frac{1 i m}{N}\left(Z / d^{N} Z\right)$, as in $\S 2$. We call a function $f(x, s)$ on $X \times U$ (where $U$ is a subset of $\Omega_{p}$ with no isolated points) 1ocally analytic if every ( $a, b$ ) $\epsilon X \times U$ has a neighborhood $\left(a+d p^{N} Z_{p}\right) \times\left(D_{b}(r) \cap U\right)$ on which $f(x, s)=$ $\sum a_{i j}(x-a)^{i}(s-b)^{j}$. Then it is easy to show that
$F(s)=\lim _{n \rightarrow \infty} \frac{1}{d p^{n}} \sum_{0 \leq j<d p} n^{f(j, s)}$
exists and is a locally analytic function of $s \in U$ (see [22]). In the case of $L_{p}(s, X)$, we had $U=Z_{p}$ and

$$
f(x, s)= \begin{cases}\left\langle x^{s} X(x)\right. & \text { if } x \in X^{*} \\ 0 & \text { otherwise }\end{cases}
$$

A special case of this construction leads to the p-adic log gamma function. We start with a function on $Z_{p} \times U$ of the form $f(x+s)$ with $U$ now a subset of $\Omega_{p}$ which is invariant under translation by $Z_{p}$ -

Lemma. Suppose that $f(x)$ is locally analytic on $s+z_{p}$ for some fixed $s \in \Omega_{p}$. Let
$F(s)=\lim _{n \rightarrow \infty} p^{-n} \sum_{0 \leq j<p} n(s+j)$.
Then $F$ is locally analytic on $s+Z_{p}$, and

$$
F(x+1)-F(x)=f^{\prime}(x)
$$

The proof is easy; the last assertion follows because $F(x+1)-F(x)=\lim \left(f\left(x+p^{n}\right)-f(x)\right) / p^{n}$.

The classical $\log$ gamma function satisfies $\log \Gamma(x+1)-$ $\log \Gamma(x)=\log x$. So, by the lemma, the natural way to obtain a $p$-adic analog is to let $f^{\prime}(x)=\ln _{p} x$, i.e., $f(x)=x \ln _{p} x-x$ (see the remark at the end of §I.3). Thus, J. Diamond [22] defined his p-adic log gamma function as

$$
\begin{equation*}
G_{p}(x)=\lim _{n \rightarrow-\infty} p^{-n} \sum_{0 \leq j<p}^{n}(x+j) \ln _{p}(x+j)-(x+j) \tag{7.4}
\end{equation*}
$$

for $x \in \Omega_{p}-Z_{p}$. Thus,
$G_{p}(x+1)-G_{p}(x)=\ln _{p} x$.
Note that it is inevitable that a continuous $p$-adic function satisfying (7.5) not be defined on $Z_{p}$. Namely, $\ln _{p} 0$, and hence either $G_{p}(1)$ or $G_{p}(0)$, is not defined. It then follows by
induction that $G_{p}$ cannot be defined either on the positive integers or the negative integers, both of which are dense in $Z_{p}$.

The other possible candidate for a p-adic log ganma function, namely $\ln _{p} \Gamma_{p}$, is defined on $Z_{p}$, but it only satisfies (7.5) when $x \in \underset{p}{p \times}$ (see property (1) $\stackrel{p}{\circ}{ }_{p} r_{p}$ ).

The two functions $G_{p}$ and $\ln _{p} P_{p}$ are related as follows. First note that $G_{p}(x)+G_{p}(1-x)=0$, as follows inmediately from (7.4) (after replacing $j$ by $p^{n}-1-j$ ). I now claim: If $x \in Z_{p}$, , then

$$
\begin{equation*}
\ln _{p} r_{p}(x)=\sum_{0 \leq i<p, p Y_{i+x}} G_{p}\left(\frac{i+x}{p}\right) \tag{7.6}
\end{equation*}
$$

i.e., we omit the one value of $i$ for which $G_{p}((i+x) / p)$ is not defined because $(i+x) / p \in Z_{p}$. To prove (7.6), we note that both sides vanish when $x=0$, since $G_{p}(i / p)+G_{p}((p-i) / p)=0$; and both sides of (7.6) change by the same amount when $x$ is replaced by $x+1$, namely by $\ln _{p} x$ if $x \in Z_{p}^{*}$ and by 0 if $x \in \mathcal{p}_{p}$. Since the nonnegative integers are dense in $Z_{p}$, we have (7.6) for all $\mathrm{x} \in \mathrm{Z}_{\mathrm{p}}$.

We discover an interesting relationship between $G_{P}$ and the zeta function if we expand $G$ in powers of $1 / x$ for $x$ large. Suppose $|x|_{p}>1$. We have

$$
\begin{aligned}
G_{p}(x)= & \lim _{n \rightarrow \infty} p^{-n} \sum_{0 \leq j<p}(x+j) \ln _{p} x+ \\
& x \lim _{n \rightarrow \infty} p^{-n} \sum_{0 \leq j<p}\left(1+\frac{j}{x}\right)\left(-1+\ln _{p}\left(1+\frac{j}{x}\right)\right) \\
= & \left(x-\frac{1}{2}\right) \ln _{p} x-x+x_{n \rightarrow-\infty} \lim _{p} \sum_{0 \leq j<p} \sum_{k=1}^{\infty}(-1)^{k+1}\left(\frac{j}{x}\right)^{k+1}\left(\frac{1}{k}-\frac{1}{k+1}\right) \\
= & \left(x-\frac{1}{2}\right) \ln _{p} x-x+\sum_{k=1}^{\infty} \frac{B_{k+1}}{k(k+1)} x^{-k} \\
= & \left(x-\frac{1}{2}\right) \ln _{p} x-x+\frac{1}{12 x}-\frac{1}{360 x^{3}}+\ldots,
\end{aligned}
$$

where we used (7.1) and the fact that $B_{k}=0$ for odd $k \geq 3$.

Hence;

$$
G_{p}(x)=\left(x-\frac{1}{2}\right) \ln _{p} x-x-\sum_{k=1}^{\infty} \zeta(-k) \frac{x^{-k}}{k} .
$$

Roughly speaking, one might expect that, since $\zeta(-\mathrm{k})$ is an integral of $t^{k}, G_{p}(x)$ is essentially $-\sum \int t^{k} \frac{x^{-k}}{k} d \mu(t)=$ $\int \ln _{p}\left(1-\frac{t}{x}\right) d \mu(t)$. We shall later look more carefulity at this possibility (see $\$ 8$ and the Appendix).

Remark. In the classical case, Stirling's formula

$$
n!=\sqrt{2 \pi n} \frac{n^{n}}{e^{n}} e^{\theta / 12 n}, \quad 0<\theta<1,
$$

## gives

$$
\log \frac{P(x)}{\sqrt{2 \pi}}=\left(x-\frac{1}{2}\right) \log x-x+\frac{\theta}{12 x} .
$$

Thus, $G$ is actually the analog of $\log (\Gamma(x) / \sqrt{2 \pi})$. (From a number theoretic point of view it is often natural to normalize the gamma function by dividing by $\sqrt{2 \pi}$. For example, $\Gamma\left(\frac{1}{2}\right) / \sqrt{2 \pi}=1 / \sqrt{2}$ is algebraic; also, the right side of the Gauss multiplication formula becomes simpler, see property (5) at the beginning of $\$_{6}$.) Note that in the classical case the series (7.7) is only an asymptotic series. We cannot simply evaluate (7.7) at $x \in C$, since it diverges for all $x$ : $\left|B_{k}\right|$ grows roughly like $k!$, in contrast to $\left|B_{k}\right|_{p}$, which is bounded.

Finally, we note the following "distribution property" of $G_{p}$ which follows immediately from the definition (7.4) and the fact that $\ln _{p} p=0$ :

$$
\begin{equation*}
G_{p}(x)=\sum_{0 \leq i<p} G_{p}\left(\frac{x+i}{p}\right) \quad \text { for } \quad x \in \Omega_{p}-Z_{p} \tag{7.8}
\end{equation*}
$$

## 8. A formula for $L_{p}^{\prime}(0, X)$

The purpose of this section is to prove a formula for $\mathrm{L}_{\mathrm{p}}^{\prime}(0, x)$ which is analogous to a classical formula of Lerch (see [97], p. 271) :

$$
\begin{aligned}
L^{\prime}(0, \chi) & =B_{I, \chi} \log d+\sum_{0<a<d} X(a) \log \Gamma(a / d) \\
& =-L(0, \chi) \log d+\sum_{0<a<d} x(a) \log \Gamma(a / d) .
\end{aligned}
$$

We start by defining a twisted version of $G_{p}$ :

$$
\begin{equation*}
G_{p, z}(x)=\lim _{n \rightarrow \infty} \frac{1}{r p^{n}} \sum_{0 \leq j<r_{p} n^{n}} z^{j}(x+j)\left(\ln _{p}(x+j)-1\right), \tag{8.1}
\end{equation*}
$$

where $z^{r}=1, x \in \Omega_{p}-Z_{p}$. In particular, $G_{p, 1}=G_{p}$. The following properties of $G_{p, z}$ are proved in the same way as the analogous properties of $G_{p}$.

Proposition. The 1imit (8.1) exists for $x \in \Omega_{p}-Z_{p}$ and satis fies:

$$
\begin{align*}
& z G_{p, z}(x+1)-G_{p, z}(x)=1 n_{p} x \text { for } x \notin z_{p} ;  \tag{8.2}\\
& G_{p, z}(x)=\sum_{i=0}^{p-1} z^{i} G_{p, z}\left(\frac{x+i}{p}\right) \text { for } x \notin Z_{p} ;  \tag{8.3}\\
& G_{p, z}(x)=B_{1, z} \ln _{p} x+\sum_{k=1}^{\infty} \frac{(-1)^{k}}{k} x^{-k} L(-k, 1, z) \tag{8.4}
\end{align*}
$$

for $\cdot|x|_{p}>1$, where $B_{1, z}=1 /(z-1)=-L(0,1, z)$. (See formula (1.4) with $k=0, f$ trivial, $d=1, z=\varepsilon$; in (8.4) we are supposing that $z \neq 1$.)

We now give an expression for $G_{p, z}$ in terns of the measure $\mu_{z}$. Here $z^{r}:=1, \quad z \neq 1$, and $\mu_{z}$ is the measure on $Z_{p}$ defined by $\mu_{z}\left(a+p^{N} Z_{p}\right)=z^{a} /\left(1-z^{p^{N}}\right)$.

## Proposition.

$$
\begin{equation*}
G_{p, z}(x)=-\int_{Z_{p}} \ln _{p}(x+t) d \mu_{z}(t) \quad \text { for } x \in \Omega_{p}-Z_{p} \tag{8.5}
\end{equation*}
$$

Proof. Let $\tilde{G}_{p, z}(x)$ denote the function on the right in (8.5). Then for $|x|_{p}>1$ we have

$$
\begin{aligned}
\tilde{G}_{p, z}(x) & =-\int_{Z_{p}} \ln _{p} x d \mu_{z}(t)-\int_{Z_{p}} \ln _{p}\left(1+\frac{t}{x}\right) d \mu_{z}(t) \\
& =-\ln _{p} x \mu_{z}\left(Z_{p}\right)+\sum_{k \geq 1} \frac{1}{k}(-1)^{k} \int_{Z_{p}} t^{k} d \mu_{z}(t) x^{-k}
\end{aligned}
$$

$=-L(0,1, z) \ln _{p} x+\sum_{k \geq 1} \frac{(-1)^{k}}{k} x^{-k} L(-k, 1, z)$
by (2.5). Thus, by (8.4), $\tilde{G}_{p, z}(x)=G_{p, z}(x)$ for $|x|_{p}>1$. Now let $U_{n}=\left\{x \in \Omega_{p}| | x-\left.j\right|_{p}>p^{-n}\right.$ for all $\left.j \in Z_{p}\right\}$. Then $\Omega_{p}-Z_{p}=\bigsqcup_{n=0}^{\infty} U_{n}$. We prove that $\tilde{G}_{p, z}(x)=G_{p, z}(x)$ for $x \in U_{n}$ by induction on $n$. We just proved this equality for $n=0$. If we show that $\tilde{G}_{p, z}$, like $G_{p, z}$, satisfies (8.3), then the induction step will follow, since $x \in U_{n+1} \Rightarrow(x+i) / p \in U_{n} \quad$ for $i=0,1, \ldots, p-1$. But the change of variables $u=p t+i$ gives

$$
z^{i} \int_{Z_{p}} \ln _{p}\left(\frac{x+i}{p}+t\right) d \mu \mu_{z^{2}}(t)=\int_{i+p Z_{p}} \ln _{p}(x+u) d \mu_{z}(u)
$$

if we use the fact that $\ln _{p} p=0$ and the definition of $\mu_{z}$ (as in (3.2)). The property (3.4) then follows for $\tilde{G}_{p, z}$, and the proposition is proved.

Remark. If we define the convolution $g$ of $f$ with $\mu$ by $g(x)=\int_{Z_{p}} f(x+t) d \mu(t)$ for $x \in \Omega_{p}$ such that $f$ is continuous on $x+z_{p}$, then it follows from (2.3) and the definition of $\mu_{z}$ that $\mathrm{zg}(\mathrm{x}+1)-\mathrm{g}(\mathrm{x})=-\mathrm{f}(\mathrm{x})$ when $\mu=\mu_{z}$. Thus, if we take the preceding proposition as the definition of $G_{p, z}$, then property (8.2) follows from this equality with $f=-\ln n_{p}$.

Corollary.
$G_{p, z}^{(k)}(x)=(-1)^{k}(k-1)!\int_{Z_{p}} \frac{d \mu_{z}(t)}{(x+t)^{k}} \quad$ for $\quad x \in \Omega_{p}-Z_{p}, \quad k \geq 1$. (8.6)
Theorem (Diamond [22] and Ferrero-Greenberg [29]). Let $X$ be $\frac{\text { a nontrivial character of conductor }}{\sum} \mathrm{d}$. Then
$L_{p}^{\prime}(0, \chi)=\sum_{0<a<p d, p \nmid a} X_{1}(a) G_{p}\left(\frac{a}{p d}\right)-L_{p}(0, \chi) 1 n_{p} d$.
Proof. In order to relate twisted and untwisted $G_{p}$, we need
a lemma. As usual, $r$ is any positive integer prime to pd.
Lemiria. Let $0<a<p d, p \| a, B_{1}(x)=x-\frac{1}{2}$, and define $a^{r}$, $0<a^{i}<p d, \quad b y: r a^{\prime} \equiv a$ (mod pd). Then

$$
\begin{equation*}
\sum_{z^{r}=1} z_{G_{G}}^{a} z^{d d}\left(\frac{a}{p d}\right)=r\left(G_{p}\left(\frac{a^{\prime}}{p d}\right)+\ln _{p} r \cdot B_{1}\left(\frac{a^{\prime}}{p d}\right)\right) . \tag{8.8}
\end{equation*}
$$

To prove (8.8), we use (8.1) to write the left side as

$$
\begin{aligned}
& \lim _{n \rightarrow \infty} \frac{1}{r p^{n}} \sum_{0 \leq j<r p} \sum_{z^{r}=1} z_{z^{a+p d j}\left(\frac{a}{p d}+j\right)\left(\ln _{p}\left(\frac{a}{p d}+j\right)-1\right)} \\
&=r \lim _{n \rightarrow \infty} \frac{1}{r_{p}{ }^{n}} \sum_{0 \leq j<r p} n\left(\frac{a}{p d}+j\right)\left(\ln _{p}\left(\frac{a}{p d}+j\right)-1\right) \\
&=\lim _{n \rightarrow-\infty} p^{-n} \sum_{0 \leq j<p} n\left(\frac{a}{p d}+a^{n}+r j\right)\left(\ln _{p}\left(\frac{a}{p d}+a^{\prime \prime}+r j\right)-1\right),
\end{aligned}
$$

where $0<a^{\prime \prime}<r, a^{\prime \prime} \equiv-a / p d(\bmod r)$. Since $\frac{1}{r}\left(\frac{a}{p d}+a^{\prime \prime}\right)=\frac{a^{\prime}}{p d}$, we find that the left side of (8.8) equals

$$
\begin{aligned}
& r \lim _{n \rightarrow \infty} p^{-n} \sum_{0 \leq j<p^{n}}\left(\frac{a^{\prime}}{p d}+j\right)\left(\ln _{p} r+\ln _{p}\left(\frac{a^{\prime}}{p d}+j\right)-1\right) \\
& =r\left(B_{1}\left(\frac{a^{\prime}}{p d}\right) \operatorname{In}_{p} r+G_{p}\left(\frac{a^{\prime}}{p d}\right)\right) .
\end{aligned}
$$

We now proceed to the proof of the theorem. First, a twisted version of (8.7) (for $z \neq 1$ ) follows immediately by differentiating under the integral sign in the definition

$$
L_{p}(s, X, z)=\int_{X^{*}}\langle x\rangle^{-s} X_{1}(x) d \mu_{z}(x)
$$

and then setting $s=0$. Nanely, we have:

$$
\begin{aligned}
L_{p}^{\prime}(0, X, z) & =-\int_{X^{*}} \ln _{p} x X_{1}(x) d \mu_{z}(x) \\
& =-\ln _{p}(p d) \int_{X^{*}} X_{1} d \mu_{z}-\sum_{0<a<p d} X_{1}(a) \int_{a+p d Z_{p}} \ln _{p}\left(\frac{x}{p d}\right) d \mu_{z}(x) \\
& =-L_{p}(0, \chi, z) n_{p} d-\sum_{\substack{0<a<p d \\
p l a}} X_{1}(a) z^{a} \int_{Z_{p}} \ln n_{p}\left(\frac{a}{p d}+t\right) d \mu_{z^{p d}}(x)
\end{aligned}
$$

$=-I_{p}(0, \chi, z) 1_{n_{p}} d+\sum_{0<a<p d, p \nmid a} X_{1}(a) z^{a} G_{p, z^{p d}}\left(\frac{a}{p d}\right)$.
Now let $A$ denote the right side of (8.7). We must show that $L_{p}^{\prime}(0, X)=A$. Summing the twisted version of (8.7) over $z \neq 1$ with $z^{r}=1,{ }^{r}$ and adding $A+L_{p}^{\prime}(0, X)$ to both sides, we obtain
$A+\sum_{z^{r}=1}^{s} L_{p}^{\prime}(0, \chi, z)=L_{p}^{\prime}(0, \chi)-\sum_{z^{r}=1} L_{p}(0, \chi, z) 1 n_{p} d$

$$
+\sum_{\substack{0<a<p d \\ p l a}} x_{1}(a) \sum_{z^{T}=1} z^{a} G_{p, z} p d\left(\frac{a}{p d}\right) .
$$

Note that the relation (4.5) between the twisted and untwisted $L_{p}$ gives

$$
\langle r\rangle^{1-s} X(r) L_{p}(s, \chi)=\sum_{z^{r}=1} L_{p}(s, \chi, z),
$$

and, if we differentiate,

$$
-\ln _{p} r\langle r\rangle^{1-s} X(r) I_{p}(s, X) \div\langle r\rangle^{1-s} X(r) I_{p}^{\prime}(s, X)
$$

$$
\begin{equation*}
=\sum_{z^{r}=1} L_{p}^{\prime}(s, X, z) \tag{8.10}
\end{equation*}
$$

Using (8.9) with $s=0$ and the 1 enma (8.8), we have

$$
\begin{aligned}
A+\sum_{z^{r}=1} L_{p}^{\prime}(0, \chi, z)= & L_{p}^{\prime}(0, \chi)-\left\langle r>\chi(r) L_{p}(0, \chi) \ln _{p} d\right. \\
& +r \sum_{\substack{0<a<p d \\
p \nmid a}} X_{1}(a)\left(G_{p}\left(\frac{a^{\prime}}{p d}\right)+\left(\frac{a^{\prime}}{p^{d}}-\frac{1}{2}\right) \ln _{p} r\right) .
\end{aligned}
$$

Note that $\left\langle r>X^{(r)}=r X_{1}(r)\right.$ and $X_{1}\left(a^{\prime}\right)=X_{1}(a / r)$. Now using (8.10) with $s=0$, we obtain

$$
\begin{aligned}
A+ & r X_{1}(r) L_{p}^{\prime}(0, \chi)-r X_{1}(r) \ln _{p} r L_{p}(0, \chi)=L_{p}^{\prime}(0, \chi)- \\
& r X_{I}(r) L_{p}(0, \chi) \ln _{p} d+r X_{1}(r) \sum_{0<a^{\prime}<p d, p \nmid a^{\prime}} X_{1}\left(a^{\prime}\right) G_{p}\left(\frac{a^{\prime}}{p d}\right) \\
& +r X_{1}(r) \ln _{p} r \sum_{0<a^{\prime}<p d, p \mid a^{\prime}} X_{1}\left(a^{\prime}\right) \frac{a^{\prime}}{p d} .
\end{aligned}
$$

Since the last sum on the right equals $B_{1, X_{I}}=-L_{p}(0, X)$, we cancel that term and obtain
$A+r X_{1}(r) L_{p}^{\prime}(0 ; X)=L_{p}^{\prime}(0, X)+r X_{1}(r) A$,
using the definition of $A$. Since $r \chi_{1}(r) \neq 1$, this gives $L_{\mathrm{p}}^{i}(0, X)=A$, and the theorem is proved.

Corollary.

$$
\begin{equation*}
L_{p}^{\prime}(0, x)=\sum_{0<a<d} X_{1}(a) \ln _{p} r_{p}\left(\frac{a}{d}\right)-L_{p}(0, x) \ln _{p} d \tag{8.11}
\end{equation*}
$$

The corollary follows immediately by using the relation (7.6) between $G_{p}$ and $\ln _{p} \Gamma_{p}$ in the formula (8.7).

Remark. In a very similar manner one can express the values of $L_{p}$ at positive integers in terms of special values of the successive derivatives of $\ln _{p} \Gamma_{p}$ (see [23], [56]). If $D^{k}$ denotes the k -th derivative, one has:
$L_{p}\left(k, X_{k-1}\right)=\frac{(-d)^{-k}}{(k-1)!} \sum_{0<a<d} X(a)\left(D^{\left.k_{1 n_{p}} P_{p}\right)\left(\frac{a}{d}\right) \quad \text { for } k \geq 1 . ~}\right.$
III. GAUSS SUMS AND THE p-ADIC GAMMA FUNCTION

1. Gauss and Jacobi sums

Let $F_{q}$ be a finite field, let $K$ be a field (such as $C$ or $\Omega_{\mathrm{p}}$, Iet
$\psi: F \underset{q}{\longrightarrow} K^{*}$
be an additive character, i.e., a nontrivial honomorphism from the additive group of $F_{q}$ to the multiplicative group $K *$, and let $X: \underset{\mathrm{q}}{\mathrm{F}} \longrightarrow \mathrm{K}^{*}$
be a multiplicative character, i.e., a homomorphism from the multiplicative group $\mathrm{F}_{\mathrm{q}}^{*}$ to K . (Warning: Characters X on $\mathrm{F}_{\mathrm{q}}^{*}$ should not be confused with the Dirichlet characters on ( $z / \mathrm{q} Z$ ) which were considered in Chapter II and were also denoted $\chi$.) The Gauss sum (in $K$ ) of $X$ and $\psi$ is defined as
$g(x, \psi)=-\sum_{x \in F_{q}^{*}} x(x) \psi(x)$.
If $X_{1}$ and $X_{2}$ are two multiplicative characters of $F_{q}$, then the Jacobi sum of $\left(X_{1}, X_{2}\right)$ is defined as
$J\left(x_{1}, x_{2}\right)=-\sum_{x \in F}, x \neq 0,1$
The Gauss and Jacobi sums are usually defined without the minus sign before the sumnation, but this definition is more convenient for our purposes.)

The Gauss and Jacobi sums satisfy the following elementary properties:
(1) If $X$ is not the trivial character, and if $\bar{X}=\chi^{-1}$ denotes the conjugate character, then
$\mathrm{g}(\chi, \psi) \mathrm{g}(\bar{\chi}, \psi)=\chi(-1) \cdot \mathrm{q}$.
(2) If $X$ is nontrivial and $K=C$, then
$|\mathrm{g}(\mathrm{X}, \psi)|=\sqrt{\mathrm{q}}$
(3) If $X_{1} \chi_{2}$ is nontrivial, then
$J\left(\chi_{1}, \chi_{2}\right)=\frac{g\left(\chi_{1}, \psi\right) g\left(\chi_{2}, \psi\right)}{g\left(\chi_{1} \chi_{2}, \psi\right)}$.
Note that the expression in property (3) does not depend on the choice of (nontrivial) additive character $\psi$.

Remark. $\mathrm{g}(\mathrm{X}, \psi)$ is the analog for $\mathrm{F}_{\mathrm{q}}$ of the gamma function on R. Namely,

$$
\Gamma(s)=\int_{0}^{\infty} x^{s} e^{-x} \frac{d x}{x},
$$

i.e., $\Gamma(s)$ is the "sum" over the (positive) multiplicative group of the field (i.e., the integral with respect to its Haar measure $d x / x)$ of the product of a multiplicative character $x \mapsto x^{s}$ and an additive character $x \mapsto \mathrm{e}^{-\mathrm{x}}$. The analog of property (1) is: $\Gamma(s) \Gamma(1-s)=\frac{\pi}{\sin (\pi s)}$, in which $\pi$ plays the role of $q$ and $\sin (\pi s)$, which is essentially $e^{\pi i s}=(-1)^{s}$, plays the role of $X(-1)$. $J\left(X_{1}, X_{2}\right)$ is the analog of the beta-function

$$
\begin{equation*}
B(r, s)=\int_{0}^{I} x^{r-1}(1-x)^{s-1} d x=\frac{\Gamma(r) \Gamma(s)}{\Gamma(r+s)} . \tag{1.2}
\end{equation*}
$$

As an illustration of this striking analogy, compare the proof of the expression (1.1) for the Jacobi sum in terms of Gauss sums with the proof of the formula (1.2) for the beta-function in terms of the gamma function. Both proofs are easy, so let's write them side-by-side:

$$
=\int_{0}^{\infty} \int_{0}^{\infty} u^{r-1} v^{s-1} e^{-u-v} d u d v \quad=\sum_{u, v \in F_{q}} x_{1}(u) X_{2}(v) \Psi(u+v)
$$

$$
=\int_{0}^{\infty} u^{r-1} e^{-u} d u \int_{0}^{\infty} v^{s-1} e^{-v} d v \quad=\sum_{u \in F_{q}} x_{1}(u) \psi(u) \sum_{v \in F_{q}} x_{2}(v) \psi(v)
$$

$=\Gamma(r) \Gamma(s)$.

$$
=g\left(x_{1}, \psi\right) g\left(x_{2}, \psi\right)
$$

The analogy between Gauss and Jacobi sums and gamma and beta functions goes deeper. The purpose of this chapter is to show that Gauss sums are essentially values of the p-adic gamma function.

## 2. Fermat curves

Let $K$ be a field containing $d$ d-th roots of unity, for example, $C$ or $\Omega_{p}$ or $F_{q}$ when $q \equiv I(\bmod d)$. We let $\mu_{d}$ denote the set of d-th roots of unity. If $K$ is of characteristic p , we mast have $\mathrm{p} \mid \mathrm{d}$. The projective Fermat curve $\mathrm{F}(\mathrm{d})$, $d>2$, is defined by $X^{d}+Y^{d}=Z^{d}$. The affine curve $F(d){ }^{\text {aff }}$ is defined by $x^{d}+y^{d}=1 \quad(x=x / Z, y=Y / Z)$. The group $\mu_{d} \times \mu_{d}$ operates on $F(d)$ and $F(d){ }^{\text {aff }}$ by

$$
\left(\xi, \xi^{\prime}\right)(x, y)=\left(\xi x, \xi^{\prime} y\right), \quad \xi, \xi^{\prime} \in \mu_{d}
$$

$$
\begin{aligned}
& \text { change of variables } \\
& u=x y, \quad v=(1-x) y \\
& x=\frac{u}{u+v}, \quad y=u+v
\end{aligned}
$$

and $H_{D R}^{1}(F(d) / Q)$, and the action of $\mu_{d} \times \mu_{d}$ on these groups. To describe the homology $H_{1}(F(d), Q)$, start with the path $\gamma_{0}:[0,1] \longrightarrow F(d){ }^{\text {aff }}, \quad \gamma_{0}(t)=\left(t, \sqrt[d]{1-t^{d}}\right)$.
Fix a primitive d-th root of unity $\xi$. Let

$$
\gamma=\gamma_{0}-(1, \xi) \gamma_{0}+(\xi, \xi) \gamma_{0}-(\xi, 1) \gamma_{0},
$$

so that $\gamma$ goes first from. ( 0,1 ) to ( 1,0 ), then from ( 1,0 ) to $(0, \xi)$, then from $(0, \xi)$ to $(\xi, 0)$, then from $(\xi, 0)$ to $(0,1)$.

Note that $F(d)$ is a nonsingular plane curve of degree $d$.
From algebraic geometry we know the following
Fact. The $2 \mathrm{~g}=(\mathrm{d}-1)(\mathrm{d}-2)$ differential forms
$\omega_{r, s}=x^{r-1} y^{s-1} \frac{d x}{y^{d-1}}, \quad 1 \leq r, s \leq d-1, \quad r+s \neq d$,

## form a basis for

$H_{D R}^{1}(F(d) / Q)=\frac{\text { differentials of the second kind }\}}{\text { \{exact differentials }\}}$,
and the g forms with $\mathrm{r}+\mathrm{s}<\mathrm{d}$ form a basis for the holomorphic forms in $H_{D R}^{1}(F(d) / Q)$.

$$
\begin{aligned}
& \text { Note that for }\left(\xi, \xi^{\mathrm{r}}\right) \in \mu_{d} \times \mu_{d} \\
& \left(\xi, \xi^{\prime}\right) * \omega_{r, s}=\xi^{r-1} \xi^{, s-1} \frac{\xi}{\xi^{\prime}{ }^{d-1}} \omega_{r, s}=\xi^{r_{r}, s} \omega_{r, s},
\end{aligned}
$$

so that $\omega_{r, s}$ is an eigen-form for $\mu_{d} \times \mu_{d}$, which acts by the character $X_{r, s} ;\left(\xi, \xi^{\prime}\right) \longmapsto \xi^{r} \xi^{\prime s}$, i.e.,

$$
\alpha^{*} \omega_{r, s}=X_{r, s}(\alpha) \omega_{r, s} \quad \text { for } \quad \alpha \in \mu_{d} \times \mu_{d}
$$

Remark. The above assertion about $\omega_{r, s}$ with $r+s<d$ forming a basis for the holomorphic differentials is true for any nonsingular plane curve of degree $d$ (see [87], p. 171-173). The assertion about differentials of the second kind can be proved as
follows
First note that the points at infinity $F(d)-F(d){ }^{\text {aff }}$ are given in the coordinates $u=Z / X=1 / x, v=Y / X=y / x$ by $u=0, v=\zeta$, where $\zeta$ runs through the $d$-th roots of -1 . Consider the (d-1)-dimensional vector space of differentials on $F(d)$ of the form $H(x, y) \frac{d x}{y^{d-1}}$, where $H(x, y)=x^{d-2} P(y / x)$ is a homogeneous polynomial of degree $d-2$. In the ( $u, v$ )-coordinates,

$$
H(x, y) \frac{d x}{y^{d-1}}=-x(x / y)^{d-1} P(y / x)\left(-\frac{d x}{x^{2}}\right)=-\frac{1}{u} P(v) \frac{d u}{v^{d-1}}
$$

which has residue $\operatorname{res}_{\zeta}=-\zeta^{1-d_{P}}(\zeta)=\zeta P(\zeta)=\Sigma_{0<i<d} a_{i-1} \zeta^{i}$ at the point at infinity $(0, \zeta)$. The sum of the residues is clearly zero, but the residues cannot all be zero unless $H(x, y)=0$, since $a_{i-1}=\frac{1}{d-1} \Sigma_{\zeta} \zeta^{-i}{ }^{\text {res }}{ }_{\zeta}$. Thus, the map from the ( $d-1$ )-dimensional vector space of $H$ 's to the ( $\mathrm{d}-1$ )-dimensional vector space of possible residues at infinity whose sum is zero, is surjective, i.e., any differential form on $F(d){ }^{\text {aff }}$ differs from a differential of the second kind (i.e., one with all residues zero) by a differential of the form $H(x, y) d x / y^{d-1}$. So for suitable $H(x, y)$, homogeneous of degree $\mathrm{d}-2$, we can write
$\omega_{r, s}-H(x, y) \frac{d x}{y^{d-I}}=a$ differential of the second kind; applying $(\xi, \xi)^{*}$ gives
$\xi^{r+s} \omega_{r, s}-H(x, y) \frac{d x}{y^{d-1}}=$ another differential of the 2nd kind. If $r+s \neq 0(\bmod d)$, we can subtract and conclude that $\omega_{r, s}$ is a differential of the second kind. Finally, to show that the 2 g differentials $\omega_{r, s}$ with $r, s<d$ are linearly independent modulo exact differentials, it suffices to use the fact that they are not exact (see (2.2) below) and they are eigen-forms for $\mu_{\mathrm{d}} \times \mu_{\mathrm{d}}$ with distinct characters.

$$
\begin{aligned}
& \text { We can thus write } \\
& \mathrm{H}_{\mathrm{DR}}^{1}(\mathrm{~F}(\mathrm{~d}) / \mathrm{Q})={ }_{1 \leq \mathrm{r}, \mathrm{~s}<\mathrm{d}, \mathrm{r}+\mathrm{s} \neq \mathrm{d}}^{\oplus} \mathrm{H}_{\mathrm{DR}}^{1}(\mathrm{~F}(\mathrm{~d}) / \mathrm{Q})
\end{aligned} \mathrm{X}_{\mathrm{r}, \mathrm{~s}} \text {, }
$$

where the space $H_{D R}^{1}(F(d) / Q){ }^{X}{ }^{X}$,s of forms on which $\mu_{d} \times \mu_{d}$ acts by $\chi_{r, s}$ is one-dimensional and is spamed by $\omega_{r, s}$.

After establishing these facts about $H_{D R}^{1}(F(d) / Q)$, it is not hard to show that the classes of the paths $\left\{\left(\xi, \xi^{\prime}\right) \gamma\right\}\left(\xi, \xi^{\prime}\right) \in \mu_{d} \times \mu_{d}$ span the homology $H_{1}(F(d), Q)$.

$$
\begin{aligned}
& \text { We now compute: } \\
& \int_{\gamma} \omega_{r, s}=\int_{\gamma_{0}} \omega_{r, s}-\int_{(1, \xi) \gamma_{0}} \omega_{r, s}+\int_{(\xi, \xi) \gamma_{0}} \omega_{r, s}-\int_{(\xi, 1) \gamma_{0}} \omega_{r, s} \\
& =\int_{\gamma_{0}}(1-(1, \xi) *+(\xi, \xi) *-(\xi, 1) *) \omega_{r, s} \\
& =\left(1-\xi^{s}+\xi^{r+s}-\xi^{r}\right) \int_{\gamma_{0}}^{\omega_{r, s}} \\
& =\left(1-\xi^{r}\right)\left(1-\xi^{s}\right) \int_{\gamma_{0}} x^{r-1} y^{s-1} \frac{d x}{y^{d-1}} . \\
& \text { But } \int_{\gamma_{0}} x^{r} y^{s-d} \frac{d x}{x}=\int_{0}^{1} \mathrm{t}^{r}\left(1-t^{d}\right)^{s / d}-1 \frac{d t}{t} \\
& =\frac{1}{d} \int_{0}^{1} u^{r / d}(1-u)^{s / d-1} \frac{d u}{u} \text { (here } u=t^{d} \text { ) } \\
& =\frac{1}{\mathrm{~d}} \mathrm{~B}\left(\frac{\mathrm{r}}{\mathrm{~d}}, \frac{\mathrm{~s}}{\mathrm{~d}}\right) \\
& \text { (see (1.2)) }
\end{aligned}
$$

Thus

$$
\begin{equation*}
\int_{\gamma}^{1 s} \omega_{r, s}=\frac{\left(1-\xi^{r}\right)\left(1-\xi^{s}\right)}{d} B\left(\frac{r}{d}, \frac{s}{d}\right) . \tag{2.2}
\end{equation*}
$$

3. L-series for algebraic varieties (not to be confused with Dirich1et L-series)

Let $V_{0}$ be a separable algebraic variety of finite type over $\mathrm{F}_{\mathrm{q}}$, and let V be obtained from $\mathrm{V}_{0}$ by extending scalars to the algebraic closure $\overline{\mathrm{F}}_{\mathrm{q}}: \quad \mathrm{V}=\mathrm{V}_{0} \otimes \overline{\mathrm{~F}}_{\mathrm{q}}$. The Frobenius map $E: \mathrm{V} \rightarrow \mathrm{V}$
is the map which raises coordinates to the q-th power (in terms of
coordinate rings, it raises variables to the q-th power and keeps coefficients fixed). Thus, the $F_{q}$-points of $V_{0}$ are the fixed points of the Frobenius $F$ in $V$.

If $g \in \operatorname{End} V$, let $|f i x(g)|$ denote the number of fixed points of $g$ on V. Thus, $\left|V_{0}\left(\mathrm{~F}_{\mathrm{q}}\right)\right|=\left|f i x\left(F^{\mathrm{n}}\right)\right|$.

Let $G$ be a finite group of automorphisms of $V_{0}$ over $F_{q}$, let $\rho: G \rightarrow G L(W)$ be a finite dimensional representation of $G$ in a vector space $W$ over a field $K$ of characteristic zero, and let $X=\operatorname{Tr}(\rho)$. Then we define

$$
L\left(V_{0} / F_{q}, G, g\right)=\exp \left(\sum_{n=1}^{\infty} \frac{T^{n}}{n} \frac{1}{|G|} \sum_{g \in G} \chi\left(g^{-1}\right)\left|f i x\left(F^{n^{n}} \circ g\right)\right|\right)
$$

$$
\in K[[T]]
$$

For example, if $G=\{1\}$, then this is merely

$$
\mathrm{Z}\left(\mathrm{v}_{0} / \mathrm{F}_{\mathrm{q}}\right)=\exp \left(\sum_{\mathrm{n}=1}^{\infty} \frac{\mathrm{T}^{\mathrm{n}}}{\mathrm{n}}\left|\mathrm{~V}_{0}\left(\mathrm{~F}_{\mathrm{q}}{ }_{\mathrm{n}}\right)\right|\right) .
$$

We shall be interested in the case of Fermat curves. Let $\mathrm{V}_{0}$ be the Zariski open subset of $F(d)$ where none of the coordinates vanish:

$$
V_{0}=\left\{(x, y) \mid x^{d}+y^{d}=1, x y \neq 0\right\} .
$$

As always, we suppose $q \equiv 1(\bmod d)$, so that $\mu_{d} \subset F_{q}$. Let $K$ be an extension of $Q_{p}$ containing $Q_{p}(\sqrt{1} \sqrt{1})$. We imbed $\mu_{d}$ in $K$ by the Teichmüller map $x \mapsto \omega(x)$. Let $G=\mu_{d} \times \mu_{d}$, which acts on $V_{0}$ by (2.1). Let $X_{1}$ and $X_{2}$ be two characters of $\mu_{d}$ with values in $K$, i.e., $X_{i}$ is of the form $x \mapsto \omega(x){ }^{a_{i}} \quad(i=1,2)$.
Let $\rho=x=x_{1} \times x_{2}$. Let $\tilde{X}_{i}(i=1,2)$ be the character on $\mathrm{F}_{\mathrm{q}}^{*}$ given by $\mathrm{x} \longmapsto \mathrm{X}_{\mathbf{i}}\left(\mathrm{X}^{(\mathrm{q}-1) / \mathrm{d}}\right)$.

Claim. The coefficient of $T$ in the exponent in (3.1) is
equal to $-J\left(\tilde{X}_{1}, \tilde{x}_{2}\right)$.
To prove this claim, note that $(x, y) \in V$ is fixed by $F \circ\left(\xi, \xi^{\prime}\right)$ whenever $\left(\xi \mathrm{x}^{\mathrm{q}}, \xi^{\prime} \mathrm{y}^{\mathrm{q}}\right)=(\mathrm{x}, \mathrm{y})$. Thus, the inner sum in (3.1) for $n=1$ is $\sum_{(x, y) \in V} x_{1}\left(x^{q-1}\right) x_{2}\left(y^{q-1}\right)$.
$x^{q-1} \epsilon \mu_{d}, y^{q-1} \epsilon \mu_{d}$
But $x^{q-1}, y^{q-1} \in \mu_{d}$ if and only if $u=x^{d}$ and $v=y^{d}$ are in $\mathrm{F}_{\mathrm{q}}^{*}$. In that case $\mathrm{X}_{1}\left(\mathrm{x}^{\mathrm{q}-1}\right) \mathrm{X}_{2}\left(\mathrm{y}^{\mathrm{q}-1}\right)=\tilde{X}_{1}(\mathrm{u}) \tilde{X}_{2}(\mathrm{v})$. Since for each $(u, v)$ there are $|G|$ pairs $(x, y)$ with $u=x^{d}, v=y^{d}$, we obtain

$$
\begin{aligned}
& \frac{1}{|G|} \sum_{g \in G} x\left(g^{-1}\right)\left|\operatorname{fix}\left(F^{\circ} \mathrm{g}\right)\right|=\sum_{\substack{u, v \in F_{\mathrm{q}}^{*} \\
u \\
u \\
v=1}} \tilde{X}_{1}(u) \tilde{X}_{2}(v) \\
& =\sum_{\substack{u \in F_{q}^{\star}, u \neq 0,1}} \tilde{X}_{1}(u) \tilde{x}_{2}(1-u) \\
& =-\mathrm{J}\left(\tilde{\mathrm{x}}_{1}, \tilde{\mathrm{x}}_{2}\right) \text {. }
\end{aligned}
$$

We shall also be interested in the case of so-called "ArtinSchreier curves". Just as the Fermat curve is connected with Jacobi sums, similarly the Artin-Schreier curve

$$
\begin{equation*}
y^{p}-y=x^{d}, \quad p / d, \tag{3.2}
\end{equation*}
$$

is connected with Gauss sums. Let $A(d, p)$ denote the complete nonsingular model of the plane curve (3.2). Note that here $p$ appears in the form of the equation. This is related to the fact that Gauss sums, unlike Jacobi sums, depend on an additive character $\psi$.

In the Gauss sums we study, the additive character $\psi$ is always assumed to be of the form $\psi: \underset{\mathrm{F}}{\mathrm{q}} \xrightarrow[\mathrm{Tr}]{\mathrm{Tr}} \underset{\mathrm{P}}{\Psi_{0}} \mathrm{~K}$ * for some character $\psi_{0}$ of $F_{p}$, i.e., it is obtained by pulling back an additive character of $F_{p}$ by means of the trace map from $F_{q}$ to $F_{p}$.

The curve $A(d, p)$ is a degree $d$ covering of the $y$-line which is totally ramified over the $p$ solutions of $y^{p}-y=0$ and over the point at infinity, and is unramified elsewhere. It follows from the Hurwitz genus formula (see, e.g., [40], p. 301): that $2 \mathrm{~g}-2=$ $-2 d+(p+1)(d-1)$, so that $2 g=(d-1)(p-1)$.

Over a field of characteristic $p$ containing $\mu_{d}$, the curve $\mathrm{A}(\mathrm{d}, \mathrm{p})$ has two types of automorphisms:

$$
\begin{align*}
& \quad x \mapsto \xi x, \quad y \mapsto y, \quad \xi \in \mu_{d} ;  \tag{3.3}\\
& \text { and } \\
& x \mapsto x, \quad y \mapsto y+\alpha, \quad \alpha \in F_{p} . \tag{3.4}
\end{align*}
$$

We shall see that the first type corresponds to the multiplicative character $X$ and the second type to the additive character $\psi$ in the Gauss sum. Note that, as in the case of $F(d)$, the number $2 \mathrm{~g}=(\mathrm{d}-\mathrm{I})(\mathrm{p}-1)$ is the number of pairs $(\chi, \psi)$ with both $X$ and $\psi=\psi_{0} \circ \operatorname{Tr} \quad$ nontrivial.

Let $V_{0}=\left\{(x, y) \mid y^{p}-y=x^{d}, \quad x \neq 0\right\}$, let $G=\mu_{d} \times Z / p Z$, and let $\rho=\chi:(\xi, \alpha) \longmapsto \chi_{1}(\xi) \psi_{0}(\alpha)$.

Claim. The coefficient of $T$ in the exponent of (3.1) is equal to $-\mathrm{g}\left(\tilde{\mathrm{X}}_{1}, \psi\right)$.

This claim is proved in a manner similar to the previous one:

$$
\begin{aligned}
\sum_{\mathrm{g} \in \mathrm{G}} \mathrm{X}\left(\mathrm{~g}^{-1}\right) \mid \operatorname{fix}\langle(F \circ g)| & =\sum_{\substack{(\xi, \alpha),(x, y) \\
y^{p}-y=x^{d} \neq 0 \\
(x, y)=\left(\xi x^{q}, y^{q}+\alpha\right)}} x^{-1}(\xi, \alpha) \\
& =\sum_{\substack{(x, y)}} x_{1}\left(u^{(q-1) / d}\right) \psi_{0}\left(y^{q}-y\right) \\
& y^{p}-y=x^{d}=u \in F_{q}^{*} \\
& =|G| \sum_{u \in F_{q}^{*}} \tilde{X}_{1}(u) \psi_{0}\left(\operatorname{Tr}_{F_{q}} / F_{p}^{u)}\right.
\end{aligned}
$$

as desired.
We now look more closely at L-series of this type. In particular, we show that the "bad" points of $F(d)$ and $A(d, p)$ that were omitted in $V_{0}$ (the points where a coordinate is zero) can in fact be included without changing the l-series, if our characters are nontrivial.

We return to the general case of a variety $V_{0}$. The Frobenius $F$ acts on the $\overrightarrow{\mathrm{F}}_{\mathrm{q}}$-points of $\mathrm{V}_{0}$. A closed point x of $\mathrm{V}_{0}$ is the same as an orbit of $F$; deg $x$ is the number of points in the orbit (equivalently, the degree over $F_{q}$ of the field containing the coordinates of the $\overline{\mathrm{F}}_{\mathrm{q}}$-points in the orbit).

Now suppose that $V_{0}$ is quasi-projective, and $G$ is a finite group of automorphisms of $\mathrm{V}_{0}$ over $\mathrm{F}_{\mathrm{q}}$ (hence commute with $F$ ). Let $X_{0}=V_{0} / G$ with the induced Frobenius endomorphism.

We first consider the case when $G$ has no fixed points, i.e., for all $g \in G$ and all $v \in V_{0}\left(\bar{F}_{q}\right)$, if $g \neq 1$, then $g v \neq v$. To every closed point $x_{0}$ in $x_{0}$ of degree $N$, we associate a conjugacy class $\operatorname{Frob}\left(x_{0}\right)$ in $G$ as follows. First choose an $x$ in the orbit $x_{0}$ and $a \quad v \in V_{0}\left(\bar{F}_{q}\right)$ lying over $x$. Then $F_{V}^{N}$ also lies over $x$, and so equals $g v$ for some unique $g \in G$. Changing our choice of $x$ in the orbit $x_{0}$ or our choice of $v$ lying over $x$ only changes $g$ by conjugation. Hence we obtain a conjugacy class Frob $\left(x_{0}\right)$ in $G$ depending only on $x_{0}$. (In our application later, $G$ will be abelian, so that $\operatorname{Frob}\left(\mathrm{x}_{0}\right)$ will be a welldefined element.)

Let $\rho$ be a representation of $G$ in a finite dimensional vector space $W$ over a field $K$ of characteristic zero, and let $\mathrm{x}=\operatorname{Tr} \rho$. Then

Claim.

$$
\mathrm{L}\left(\mathrm{~V}_{0} / \mathrm{F}_{\mathrm{q}}, G, \rho\right)=\Pi \frac{1}{\operatorname{Det}\left(1-\rho\left(\operatorname{Frob}\left(\mathrm{x}_{0}\right)\right) \cdot T^{\operatorname{deg} x_{0}}\right)},
$$

where the product is taken over all closed points $x_{0}$ of $X_{0}$.

To prove this claim, we take the $\log$ of both sides and use the fact that $\log \operatorname{Det}(1-M T)=-\sum \frac{\mathrm{T}^{\mathrm{n}}}{\mathrm{n}} \operatorname{Trace}\left(\mathrm{M}^{\mathrm{n}}\right) \quad$ for any matrix $M$. Then we need only show that

$$
\left.\left.\frac{1}{|\mathrm{G}|} \sum_{\mathrm{g} \in \mathrm{G}} \mathrm{X}\left(\mathrm{~g}^{-\mathrm{l}}\right) \right\rvert\, \mathrm{fix}^{\left(F_{\mathrm{g}}^{\mathrm{n}}\right.} \text { on } \mathrm{V}\right) \mid
$$

is equal to
$n$. coefficient of $T^{n}$ in $\sum_{x_{0}} \sum_{r=1}^{\infty} \frac{T^{r \operatorname{deg} x_{0}}}{r}$ Trace $\rho\left(\operatorname{Frob}\left(x_{0}\right)\right)^{r}$, which can be written

$$
\sum_{x_{0}} \text { of degree }\left.\right|_{\mathrm{n}} \mathrm{~s} \times\left(\operatorname{Frob}^{\mathrm{n} / \mathrm{s}}\left(\mathrm{x}_{0}\right)\right) \text {, }
$$

and this equality follows easily by writing the first sum as the sum of $X(g)$ over ail $g$ and $v$ with $F^{n} v=g v$ and then using the definition of $\operatorname{Frob}\left(\mathrm{x}_{0}\right)$.

Now let us allow $G$ to have fixed points. For $v \in V_{0}\left(\bar{F}_{q}\right)$,
let $I_{v}=\{g \in G \mid g v=v\}$, the "inertia group" of $v$. Let
$W^{I} v=\left\{w \in W \mid \rho(g) w=W\right.$ for all $\left.g \in I_{v}\right\}$.
Now Frob $\left(x_{0}\right)$ is only defined up to multiplication by elements of $I_{v}$, as well as conjugation; nevertheless, the determinant of $1-T^{\text {deg } x_{0}} \cdot \rho\left(\operatorname{Frob}\left(x_{0}\right)\right)$ acting on $W^{I}{ }^{\mathrm{v}}$ still depends only on $x_{0}$; and it is not hard to show that

$$
\begin{equation*}
\mathrm{I}\left(\mathrm{~V}_{0} / \mathrm{F}_{\mathrm{q}}, \mathrm{G}, \rho\right)=\prod_{x_{0}} \operatorname{Det}\left(1-\left.T^{\operatorname{deg} x_{0}} \rho\left(\operatorname{Frob}\left(x_{0}\right)\right)\right|_{W} I_{v}\right)^{-1} \tag{3.5}
\end{equation*}
$$

In our example $V_{0}=F(d), G=\mu_{d} \times \mu_{d}$, the fixed points occur
(1) when $Y=0$, in which case $I_{v}=1 \times \mu_{d}$;
(2) when $X=0$, in which case $I_{v}=\mu_{d} \times 1$;
(3) when $Z=0$, in which case $I_{v}$ is the diagonal in $\mu_{d} \times \mu_{d}$.

If our character $\rho=X_{X}=X_{r} \times X_{s}$ has the property that $X_{r}, X_{s}$, and $X_{r} X_{s}$ are all nontrivial (i.e., $0<r, s<d, r+s \neq d$ ), then in all cases $I_{v}$ acts nontrivially on the one-dimensional space $W$. Hence $W^{I} v=0$, and there are no contributions to (3.5) from the points with zero $X, Y$ or $Z$ coordinate.

Similarly, when $V_{0}=A(d, p), G=\mu_{d} \times z / p Z$, the fixed points are
(1) the point at infinity, where $I_{v}=G$;
(2) the points with $x=0$, where $I_{v}=\mu_{d} \times\{0\}$.

Again there is no contribution to (3.5) when the characters are nontrivial.

We now return to the general case of a variety $V_{0}$. Let $\rho_{\text {reg }}$ be the regular representation of $G$. We have

$$
\text { Trace } \rho_{\text {reg }}(g)=\left\{\begin{array}{ccc}
|G| & \text { if } & g=1 \\
0 & \text { if } & g \neq 1
\end{array}\right.
$$

It is immediate from the definition that

$$
\mathrm{L}\left(\mathrm{v}_{0} / \mathrm{F}_{\mathrm{q}}, \mathrm{G}, \rho_{\mathrm{reg}}\right)=\mathrm{Z}\left(\mathrm{~V}_{0} / \mathrm{F}_{\mathrm{q}}\right)=\exp \sum \frac{\mathrm{T}^{\mathrm{n}}}{\mathrm{n}}\left|\mathrm{v}_{0}\left(\mathrm{~F}_{\mathrm{q}}\right)\right| ;
$$

and aiso

$$
L\left(V_{0} / F_{q}, G, \rho_{\text {trivial }}\right)=Z\left(X_{0} / F_{q}\right) \quad\left(\text { recall } \quad x_{0}=V_{0} / G\right)
$$

Since trivialiy we have

$$
L\left(V_{0} / F_{\mathrm{q}}, G, \rho_{1} \oplus \rho_{2}\right)=L\left(V_{0} / \mathrm{F}_{\mathrm{q}}, \mathrm{G}, \rho_{1}\right) \cdot \mathrm{L}\left(\mathrm{~V}_{0} / \mathrm{F}_{\mathrm{q}}, \mathrm{G}, \rho_{2}\right),
$$

it follows that the decomposition

$$
\rho_{\text {reg }}=\oplus \rho^{\operatorname{deg} \rho},
$$

where the summation is over all irreducible representations of $G$, gives a corresponding product decomposition of $Z\left(V_{0} / F_{q}\right)$ :

$$
\begin{equation*}
z\left(V_{0} / F_{q}\right)=Z\left(X_{0} / F_{q}\right)\left\lceil L\left(v_{0} / F_{q}, G, \rho\right)^{\operatorname{deg}} \rho,\right. \tag{3.6}
\end{equation*}
$$

where the product is over all nontrivial irreducible representations
$\rho$.
Suppose $\mathrm{V}_{0}$ is a projective, nonsingular, geometrically connected curve over $F_{q}$, of genus $g$. Then $X_{0}$, the quotient of $V_{0}$ by $G$, is also nonsingular, say of genus $g^{\prime}$. We have

$$
\mathrm{Z}\left(\mathrm{~V}_{0} / \mathrm{F}_{\mathrm{q}}\right)=\frac{\text { polynomial in } \mathrm{Z}[\mathrm{~T}] \text { of degree } 2 \mathrm{~g}}{(1-T)(1-\mathrm{qT})}
$$

Moreover, if $\rho$ is irreducible and nontrivial, then

$$
L\left(V_{0} / F_{q}, G, P\right) \text { is a polynomial in } T .
$$

This was proved for curves by Weil in the $1940^{\prime} \mathrm{s}$, but it is only a conjecture in the general case of higher dimensional $\mathrm{V}_{0}$. It then follows that

$$
\sum \operatorname{deg} \rho \operatorname{deg} L\left(V_{0} / F_{q}, G, \rho\right)=2 g-2 g^{\prime},
$$

where the summation is over all irreducible nontrivial representations $\rho$.
4. Cohomology

Let $Z_{p}^{\text {unr }}=0_{Q_{p}^{\text {unr }}}=\left\{\left.x \in Q_{p}^{\text {unr }}| | x\right|_{p} \leq 1\right\}$; thus, $Z_{p}^{\text {unr }}$ is the ring extension of $Z_{p}$ generated by all $N$-th roots of unity with $\mathrm{p} \mid \mathrm{N}$.

Fact. For every prime \& there exists a functor $H^{1}$

$$
\left.\begin{array}{l}
\text { projective nonsingular } \\
\begin{array}{l}
\text { geometrically connected } \\
\text { curves of genus } \mathrm{g} \\
\text { over }
\end{array} \overline{\mathrm{F}}_{\mathrm{q}}
\end{array}\right\} \longrightarrow\left\{\begin{array}{l}
\text { free modules (of rank } 2 \mathrm{~g} \text { ) } \\
\text { over } \begin{cases}\mathrm{z}_{\ell} \text { if } \ell \neq \mathrm{p} \\
\mathrm{z}_{\mathrm{p}}^{\text {unr }} & \text { if } \ell=\mathrm{p}\end{cases}
\end{array}\right.
$$

(namely, $H_{\text {étale }}^{1}\left(V, Z_{\ell}\right)$ for $\ell \neq p$ and $H_{\text {crystalline }}^{\left(V / Z_{p}^{u n r}\right)}$ for $\ell=p$ ), such that if $V=V_{0}\left(\otimes_{F_{q}} \bar{F}_{q}\right.$ with $V_{0}$ defined over $F_{q}$, and if $F$ is the $q$-th power Frobenius endomorphism, then

$$
\operatorname{Trace}\left(F^{*} \mid \mathrm{H}^{1}(\mathrm{~V})\right)=\mathrm{I}+\mathrm{q}-\left|\mathrm{V}_{0}\left(\mathrm{~F}_{\mathrm{q}}\right)\right|
$$

This fact implies that

$$
\begin{aligned}
\mathrm{Z}\left(\mathrm{~V}_{0} / \mathrm{F}_{\mathrm{q}}\right) & =\exp \sum_{\mathrm{n}=1}^{\infty} \frac{\mathrm{T}^{n}}{\mathrm{n}}\left|\mathrm{~V}_{0}\left(\mathrm{~F}_{\mathrm{q}^{n}}\right)\right| \\
& =\exp \sum_{\mathrm{n}=1}^{\infty} \frac{T^{n}}{\mathrm{n}}\left(1+\mathrm{q}^{\mathrm{n}}-\operatorname{Trace}\left(F *^{n} \mid H^{1}(\mathrm{~V})\right)\right) \\
& =\frac{1}{1-T} \cdot \frac{1}{1-\mathrm{qT}} \operatorname{Det}\left(1-T \cdot F * \mid H^{I}(V)\right) .
\end{aligned}
$$

Here the Det term is a polynomial which is clearly in $\mathrm{Z}[\mathrm{T}]$.

$$
\text { Similarly, for L-series one can construct for } g \in G \text { a "twisted" }
$$ variety $V_{0}^{\prime}$ defined over $F_{q}$ such that $V_{0}^{\prime}{ }^{\otimes}{\underset{F}{q}}^{F_{q}} \overline{\mathrm{~F}}^{\sim} \simeq V_{0}{ }^{\otimes}{\underset{F}{F}} \overline{\mathrm{~F}}_{\mathrm{q}}=\mathrm{V}$, while the Frobenius for $V_{0}^{\prime}$ is $F \circ g \quad(F$ is the Frobenius for $\mathrm{V}_{0}$ ). Thus, (4.1) implies that

$$
\begin{equation*}
\operatorname{Trace}\left(\left(F^{\mathrm{n}} \circ \mathrm{~g}\right) * \mid \mathrm{H}^{1}(\mathrm{~V})\right)=1+\mathrm{q}^{\mathrm{n}}-\left|\mathrm{fix}\left(F^{\mathrm{n}} \circ \mathrm{~g}\right)\right| \tag{4.2}
\end{equation*}
$$

Now let $\rho$ be an absolutely irreducible representation of $G$ in a vector space over a field $K$ which we assume contains $Z_{\ell}$ or $z_{p}^{\text {unr }}$. Then the subspace of $H^{1}(V) \otimes K$ on which $G$ acts by $p$ is

$$
\left\langle\mathrm{H}^{1}(\mathrm{~V}) \otimes \mathrm{K}\right)^{\rho}=\left(\frac{1}{|\mathrm{G}|} \sum_{\mathrm{g} \in \mathrm{G}} \mathrm{X}\left(\mathrm{~g}^{-1}\right) \mathrm{g}\right)\left(\mathrm{H}^{1}(\mathrm{~V}) \otimes \mathrm{K}\right),
$$

and, using (4.2) and (3.1), one easily shows that

$$
\begin{equation*}
L\left(V_{0} / F_{q}, G, \rho\right)=\operatorname{Det}\left(I-T \cdot F * \mid\left(H^{1}(V) \otimes R\right)^{\rho}\right) \tag{4.3}
\end{equation*}
$$

We now apply (3.6) and (4.3) to our examples $F(d)$ and $A(d, p)$.
It follows from (4.3) that the coefficient of $T$ in
$\mathrm{L}\left(\mathrm{V}_{0} / \mathrm{F}_{\mathrm{q}}, \mathrm{G}, \mathrm{\rho}\right)$ is - Trace $F^{*} \mid\left(\mathrm{H}^{1}(\mathrm{~V}) \otimes \mathrm{K}\right)^{\rho}$. In the case $\mathrm{V}_{0}=\mathrm{F}(\mathrm{d})$, $G=\mu_{d} \times \mu_{d}, \rho=\chi_{r} \times \chi_{s}$ with $0<r, s<d, r+s \neq d$, we showed that this coefficient is $-J\left(\tilde{X}_{r}, \tilde{X}_{s}\right) \neq 0$. Since the 2 g spaces $\left(\mathrm{H}^{1}(\mathrm{~F}(\mathrm{~d})) \otimes \mathrm{K}\right)^{\chi_{r} \times \chi_{S}}$ are nonzero, it follows that each of them is one-dimensional, and we have the direct sum decomposition
$\mathrm{H}^{1}(\mathrm{~F}(\mathrm{~d})) \oplus \mathrm{K}=\underset{\mathrm{r}, \mathrm{s}}{\oplus}\left(\mathrm{H}^{1}(\mathrm{~F}(\mathrm{~d})) \oplus \mathrm{R}\right)^{\mathrm{X}_{\mathrm{r}}{ }^{\mathrm{X}} \mathrm{X}_{\mathrm{s}}}$.
Moreover, the eigen-value of $F^{*}$ on $\left(H^{1}(F(d)) \otimes K\right)^{X_{r} \times X_{s}}$ is. Trace $F^{*}=J\left(\tilde{X}_{r}, \tilde{X}_{s}\right)$. Then (3.6) (with $V_{0}=F(d)$ and $X_{0}=V_{0} / G$ $=$ the projective line) gives

$$
Z\left(F(d) / F_{q}\right)=\frac{1}{(1-T)(1-q T)} \prod_{\substack{1 \leq r, s<d \\ r+s \neq d}}\left(1-J\left(\tilde{X}_{r}, \tilde{X}_{s}\right) T\right),
$$

as Weil explained in his famous paper in 1949 [99].

$$
\text { Next, 1et } V_{0}=A(d, p): y^{p}-y=x^{d} \text { over } F_{q}, q \equiv 1(\bmod d) ;
$$ $G=\mu_{d} \times Z / p Z$. We have shown that each of the $2 g=(d-1)(p-1)$ characters $\rho=\chi \times \psi_{0}$ with $X$ and $\psi_{0}$ nontrivial gives nonzero Trace $F^{*} \mid\left(\mathrm{H}^{1}(\mathrm{~A}(\mathrm{~d}, \mathrm{p})) \otimes \mathrm{K}\right)^{\chi \times \psi_{0}}=\mathrm{g}\left(\tilde{\mathrm{X}}, \psi_{0}{ }^{\circ} \mathrm{Tr}\right)$. Thus, we again have the decomposition

$$
\mathrm{H}^{\mathrm{I}}(\mathrm{~A}(\mathrm{~d}, \mathrm{p})) \otimes \mathrm{K}=\mathrm{X}_{\mathrm{X}} \psi_{0}^{\oplus} \text { nontrivial }\left(\mathrm{H}^{1}(\mathrm{~A}(\mathrm{~d}, \mathrm{p})) \otimes \mathrm{K}\right)^{\mathrm{X} \times \psi_{0}}
$$

into one-dimensional eigenspaces with $F$ acting by $g\left(\tilde{X}, \psi_{0} \circ \mathrm{Tr}\right)$ on the $\chi \times \psi_{0}$-component. We conclude:

$$
\mathrm{Z}\left(\mathrm{~A}(\mathrm{~d}, \mathrm{p}) / \mathrm{F}_{\mathrm{q}}\right)=\frac{1}{(1-\mathrm{T})(1-\mathrm{q} \mathrm{~T})} \quad \mathrm{X}, \psi_{0} \prod_{\text {nontrivial }}^{\left(1-\mathrm{g}\left(\tilde{\mathrm{X}}, \psi_{0}{ }^{\circ} \mathrm{Tr}\right) \mathrm{T}\right)} .
$$

We obtain a number theoretic corollary if we replace $F_{q}$ by $\mathrm{F}_{\mathrm{q}^{\mathrm{n}}}$ and replace $\mathrm{g}(\tilde{\mathrm{X}}, \psi)$ by the corresponding Gauss sum over $\mathrm{F}_{\mathrm{q}^{\mathrm{n}}}$ :

Namely, we have

$$
\mathrm{g}_{\mathrm{q}^{\mathrm{n}}}\left(\tilde{\mathrm{X}}, \psi_{0}{ }^{\circ} \mathrm{Tr}_{\left.\mathrm{F}_{\mathrm{q}} / \mathrm{F}_{\mathrm{p}}\right)} \quad=\begin{array}{l}
\text { action of the } \mathrm{q}^{\mathrm{n}} \text {-power Frobenius } \\
\text { on }\left(\mathrm{H}^{1}(\mathrm{~A}(\mathrm{~d}, \mathrm{p})) \otimes \mathrm{K}\right)
\end{array}\right.
$$

$=F *^{\mathrm{n}} \mid\left(\mathrm{H}^{1}(\mathrm{~A}(\mathrm{~d}, \mathrm{p})) \otimes \mathrm{K}\right)^{X \times \psi_{0}}$
$=\left(\mathrm{g}_{\mathrm{F}_{\mathrm{q}}}\left(\tilde{\mathrm{x}}, \psi_{0}{ }^{\circ} \mathrm{Tr}\right)\right)^{\mathrm{n}}$,
which is known as the Hasse-Davenport relation for Gauss sums.

## 5. p-adic cohomology

By explicitly constructing a p-adic $H^{I}$ for $A(d, p)$, we shall derive a p-adic expression for $g\left(\tilde{X}, \psi_{0}^{\circ}{ }^{\circ} r\right)=F * \mid\left(H^{1} \otimes K\right){ }^{\chi \times \psi_{0}}$, showing how special values of the p-adic gamma function arise as eigen-values of Frobenius. For simplicity, we shall assume $p>2$.

A key role will be played by the function $e^{\pi\left(x-x^{p}\right)}$ considered in §I.3. As before, we denote $E_{\pi}(x)=e^{\pi\left(x-x^{P}\right)}$.

Proposition. There is a one-to-one correspondence between ( $\mathrm{p}-1$ )-th roots $\pi$ of -p and nontrivial additive characters $\psi_{0}$ of $F_{p}$ such that
$\psi_{0}(1)=E_{\pi}(1) \equiv 1+\pi\left(\bmod \pi^{2}\right)$.
We then have
$\psi_{0}(a)=E_{\pi}(\omega(a)) \quad$ for $\quad a \in F_{p}$.
(Recall that $\omega(a)$ denotes the Teichmüller representative.)

$$
\text { Proof. For } x \in D(1) \text { we have: } E_{\pi}(x)^{p}=e^{p \pi\left(x-x^{p}\right)}=
$$ $\sum \frac{(p \pi)^{i}}{i!}\left(x-x^{p}\right)^{i}$. (Thus, because of the $p$ in the exponent, we can evaluate $\mathrm{E}_{\pi}(\mathrm{x})^{\mathrm{p}}$ by first evaluating the exponent and then expanding.) Hence, if $x \in D(1)$ satisfies $x-x^{P}=0 \quad-$ in other words, if $x=\omega(a)$ for some $a \in F_{p}$-- then $E_{\pi}(x)^{p}=1$. Thus, each $E_{\pi}(\omega(a))$ is a $p$-th root of 1 . It follows from the expansion $E_{\pi}(x)=1+\pi x+\frac{1}{2} \pi^{2} x^{2}+\ldots$ that $E_{\pi}(\omega(a)) \equiv 1+\pi a$ (mod $\pi^{2}$ ). The proposition now follows easily.

One similarly shows that if $\psi=\psi_{0} 0{ }^{\circ} \mathrm{Tr}_{\mathrm{F}_{\mathrm{q}}} / \mathrm{F}_{\mathrm{p}}$ is the additive character on $\mathrm{F}_{\mathrm{q}}$ and if $\beta \in \mathrm{F}_{\mathrm{q}}$, then

$$
\psi(\beta)=\prod_{0 \leq i<f} E_{\pi}\left(\omega\left(\beta^{\mathrm{P}^{i}}\right)\right),
$$

where $q=p^{f}$ and $\pi$ is the ( $p-1$ )-th root of $-p$ corresponding to $\psi_{0}$.

Let $V$ be a nonsingular algebraic variety over a perfect field $k$ of characteristic $p$. In our application. $V$ will be the curve $A(d, p)$ and $k$ will be $F_{q}$. Let $R$ be a complete discrete valuation ring with maximal ideal $M_{R}$ and residue field $k$, and let $K$ be its fraction field. For us, $R$ will be the ring of integers in $Q_{p}\left(\frac{\mathrm{q}-1}{\sqrt{1}}, \pi\right)$, where $\pi^{\mathrm{p}-1}=-\mathrm{p}$.

Washnitzer and Monsky [75] constructed an explicit version of p-adic $H^{i}$ which to a so-called "special affine open set" associates a $K$-vector space. If $U=\left\{U_{j}\right\}$ is a covering of $V$ by special affine open sets, then the map $U_{j} \mapsto H^{i}\left(U_{j}\right)$ defines a Zariski presheaf $H^{i}$ of $K$-vector spaces on $V$, such that the cohomology $\quad E_{2}^{p}, q=H^{p}\left(V, H^{q}\right) \Longrightarrow H^{\star}(V) \quad$ conjecturally abuts to a "good" p-adic cohomology (in the sense of $\S 4$ ). This has been proved in the case when $V$ is a curve. We shall only need the Washnitzer-Monsky $H^{1}$ for curves.

Theorem (Washnitzer-Monsky). The functor

$$
V \longmapsto H^{1}(V) \mathrm{de⿹}_{\mathrm{f}} \mathrm{H}^{0}\left(\mathrm{~V}, \mathrm{H}^{1}\right) \quad \text { (global sections) }
$$

on complete nonsingular geometrically connected curves V satisfies the properties:
rank $H^{1}(V)=2$ genus(V);
$\operatorname{Trace}\left(F * H^{1}(\mathrm{~V})\right)=\mathrm{q}+1-|\mathrm{fix}(F)|$,
when $\mathrm{V}=\mathrm{V}_{0} \otimes_{\mathrm{F}_{\mathrm{q}}} \mathrm{k}$ for some $\mathrm{V}_{0}$ defined over $\mathrm{F}_{\mathrm{q}}$, where $F=F_{\mathrm{F}_{\mathrm{q}}}$
is the $q$-th power Frobenius endomorphism of $V$.
The basic type of ring which goes into the construction of $H^{i}$ is $R \ll x_{1}, \ldots, x_{N} \gg$, which is defined as

$$
\begin{aligned}
\left\{\sum A_{w} x^{w} \in R\left[\left[x_{1}, \ldots, x_{N}\right]\right] \mid\right. & \text { for some real number } B \text { and } \\
& \left.\varepsilon>0 \text { we have ord } A_{w} \geq \varepsilon|w|+B\right\}
\end{aligned}
$$

where $w=\left(w_{1}, \ldots, w_{N}\right), x^{w}=x_{1}{ }^{w_{1}} \cdots x_{N}{ }_{N}, \quad|w|=\sum w_{i} . \quad$ Equivalently,

$$
\begin{aligned}
R \ll x_{1}, \ldots, x_{N} \gg
\end{aligned}=\left\{f \in R\left[\left[x_{1}, \ldots, x_{N}\right]\right] \mid \underset{\text { for }}{\exists r>1,} x_{1}, \ldots, x_{N} \in D(r)\right\} .
$$

It turns out that $H^{1}$ for the p-adic "affine line" will be $R<\langle x\rangle>\left(\otimes K / \frac{d}{d x} R<\langle x\rangle>(Q) K=0\right.$, as expected.

Remark. In p-adic analysis $D(1)$ can often be thought of as the affine line, since it is the smallest disc containing the Teichmüller representatives of the points of the affine line over $\overline{\mathrm{F}}_{\mathrm{p}}$. One might ask why the simpler ring $A=\{\mathrm{f} \in \mathrm{K}[[\mathrm{x}]\} \mid \mathrm{f}$ converges on $D(1)\}$ cannot be used. In fact, $d / d x$ is not surjective on A. (For example, $a=\sum_{p^{j}} x^{p^{j}-1}=\frac{d}{d x} \sum_{x^{p}} p^{j} \in A$, but $\sum x^{p^{j}} \frac{1}{\&} A$, so $a$ is not $\frac{d}{d x}$ of any element of A.) In fact, $\operatorname{rank}\left(\dot{A} / \frac{d}{d x} A\right)=\infty \quad$ rather than $2 \cdot$ genus(1ine) $=0$. However, $\frac{d}{d x}$ is surjective on $R<\langle x\rangle>\otimes K$.

For simplicity, suppose that the nonsingular curve $\mathrm{V}^{\text {aff }}$ is given by one equation $f(x, y)=0$ (as will be the case in our application to $\mathrm{A}(\mathrm{d}, \mathrm{p})$ ). Consider the coordinate ring of the Zaxiski open subset $U$ of $V^{\text {aff }}$ where the tangent to the curve is not vertical:

$$
k[x, y, t] /\left(f(x, y), t \frac{\partial f}{\partial y}-1\right)
$$

Such an affine open set is called a "special affine" open set $U$ over $k$. We now define the "dagger ring" for $U$ to be the quotient.
$A^{\dagger}(U)=R \ll x, y, t \gg /\left(F(x, y), t \frac{\partial F}{\partial y}-1\right)$,
where $F(x, y)$ is any fixed polynomial in $R[x, y]$ whose reduction modulo $M_{R}$ is $f(x, y)$. Note that $A^{\dagger}(U) / M_{R} A^{\dagger}(U)=k[x, y, t] /$ (f, $t \frac{\partial f}{\partial y}-1$ ), i.e., $A^{\dagger}(U)$ "lifts" the coordinate ring of the special affine open set $U$. (Of course, $A^{+}(U)$ is not unique, since $F(x, y)$ is not unique.)

In our example $V=A(d, p)$, where $f(x, y)=y^{p}-y-x^{d}$, we have $\frac{\partial f}{\partial y}=-1$. We can take $F(x, y)=y^{p}-y-x^{d}$, so that $\partial F / \partial y=$ $\mathrm{py}^{\mathrm{p}-1}-1$, which is invertible in the ring $\mathrm{R} \ll \mathrm{y} \gg$ (since $\left.\frac{1}{1-p y^{p-1}}=\sum p^{j} y^{(p-1) j}\right)$. Thus, for $U=A(d, p)^{\text {aff }}=A(d, p)-$ \{point at infinity\}, we have

$$
A^{\dagger}(U)=R \ll x, y \gg /\left(y^{P}-y-x^{d}\right) .
$$

Because $\mathrm{U}=\mathrm{V}-\{$ point $\}$, we are in an especially convenient situation for computing $H^{1}(V)$. Namely, if $V$ is a complete, nonsingular geometrically connected curve over $k$ such that $V-$ \{point\} is a special affine $U$, then

$$
H^{I}(V) \simeq H^{1}(V-\{\text { point }\})
$$

Roughly speaking, this is because, if $v_{1}, \ldots, v_{s}$ are finitely many points of $V$, then

$$
H^{1}(V) \longrightarrow H^{1}\left(V-\left\{v_{i}\right\}\right)
$$

can be identified as the subset whose residue at each point vanishes. Since the sum of the residues vanishes in $H^{1}\left(V-\left\{v_{i}\right\}\right)$; this subset is all of $H^{1}\left(V-\left\{v_{i}\right\}\right)$ if $s=1$.

We now state a general fact (see [75]) about lifting a k -morphism $\phi_{0}: \mathrm{U}_{1} \longrightarrow \mathrm{U}_{2}$ of special affine open sets to a morphism $\phi: A^{\dagger}\left(U_{2}\right) \longrightarrow A^{\dagger}\left(U_{1}\right)$ of their dagger rings. For simplicity of
notation, we suppose that there is only one x variable, as will be the case in our application. Suppose that

$$
\mathrm{t}_{i}=\operatorname{Spec} k[\mathrm{x}, \mathrm{y}, \mathrm{t}] /\left(\mathrm{f}_{\mathrm{i}}(\mathrm{x}, \mathrm{y}), \mathrm{t} \frac{\partial \mathrm{f}_{1}}{\partial \mathrm{y}}-1\right), \quad i=1,2
$$

First of all, such a lifting $\phi$ exists ( $\phi$ is not unique). Next, suppose that we fix an element $\phi(x) \in A^{\dagger}\left(U_{1}\right)$ which reduces to $\phi_{0}(\mathrm{x})$ modulo $\mathrm{M}_{\mathrm{R}}$. Then there exists a unique choice of $\phi(\mathrm{y})$ and $\phi(\mathrm{t})$ iifting $\phi_{0}(\mathrm{y}), \phi_{0}(\mathrm{t})$. Finally, let
$\Omega_{A^{+}(U)}^{\bullet}=A^{+}(U) \oplus A^{+}(U) d x$
with differential $d$, and let $H^{1}(U)=H_{A^{1}(U)}^{(X) K}=$ $A^{+}(U) \otimes K d x / d\left(A^{+}(U) \otimes K\right)$. Then $H^{1}\left(U_{i}\right)$ is independent of the choice of $F_{i}$ lifting $f_{i}$, and $\phi *: H^{1}\left(U_{2}\right) \longrightarrow H^{I}\left(U_{1}\right)$ is Endependent of the choice of $\phi$ lifting $\phi_{0}$. Thus, we may choose any convenient lifting.

We show how to choose a convenient lifting of $\phi_{0}$ in the case $\mathrm{U}=\mathrm{U}_{1}=\mathrm{U}_{2}=\mathrm{A}(\mathrm{d}, \mathrm{p})^{\text {aff }}$ and $\phi_{0}: \mathrm{x} \mapsto \mathrm{x}, \quad \mathrm{y} \mapsto \mathrm{y}+\alpha, \quad \alpha$, $\{0,1, \ldots, p-1\}$. We want to construct an endomorphism $\phi$ of $A^{+}(U)=R \ll x, y \gg /\left(y^{p}-y-x^{d}\right)$
which reduces to $\phi_{0}$ modulo $M_{R}$. It is simplest to choose $\phi(x)$
$=\mathrm{x}$. Then the above fact asserts that there exists a unique
$\phi(y) \in A^{\dagger}(U)$ such that $\phi(y) \equiv y+\alpha \quad \bmod M_{R}$. To see this concretely, we note that $z=\phi(y)-y-\alpha$ must satisfy

$$
(z+y+\alpha)^{P}-(z+y+\alpha)=y^{p}-y
$$

in other words,

$$
\begin{aligned}
z^{p}+\sum_{i=1}^{p-2}\binom{p}{i}(y+\alpha)^{i} z^{p-i}+ & \left(p(y+\alpha)^{p-1}-1\right) z \\
& +(y+\alpha)^{p}-y^{p}-\alpha=0
\end{aligned}
$$

See example (2) in §I.4.b, where we saw that if ord ${ }_{p}>-\frac{1}{p-1}$, so that $\lambda=\operatorname{ord}_{p}\left((y+\alpha)^{p}-y^{p}-\alpha\right)>0$, then there is a unique solution $z$ with $0<\operatorname{ord}_{p} z=\lambda$. This $z$ can be expressed as a power series in $y$ by first solving (5.1) mod $p^{2 \lambda}$;

$$
\left(p(y+\alpha)^{p-1}-1\right) z+(y+\alpha)^{p}-y^{p}-\alpha=0
$$

then substituting the approximate solution in place of all higher powers $z^{2}, \ldots, z^{p}$ in (5.1) and again solving the resulting linear equation for $z$, and so on. The result is a power series in $y$ with coefficients in $z_{p} \subset R$ which converges for ord $y>-\frac{1}{p-1}$. Hence, $\phi(y)=z+y+\alpha \in \quad R \ll y \gg \subset A^{\dagger}(U) \quad$ is the desired lifting of $\phi_{0}(y)=y+\alpha$.

Of course, the other type of automorphism $\phi_{0}: x \longmapsto \xi x$, $\mathrm{y} \longmapsto \mathrm{y}$, where $\xi^{\mathrm{d}}=1, \quad \xi \in \mathrm{~F}_{\mathrm{q}}$, can be lifted simply to $\phi: x \longmapsto \omega(\xi) x, y \longrightarrow y$, where $\omega(\xi) \in R$ is the Teichmïller representative.

Thus, the group $G=\mu_{d} \times Z / p Z$ acts on $A^{\dagger}(U)$, and hence on the dagger cohomology
$H^{1}(A(d, p))=A^{\dagger}(U) \otimes K d x / d\left(A^{+}(U) \otimes K\right)$.
6. p-adic formula for Gauss sums

Because we claimed that the dagger $H^{1}$ is a good cohomology, the subspace of $H^{1}(A(d, p))$ on which $G=\mu_{d} \times Z / p Z$ acts by $X \times \psi_{0}$ should be one-dimensional (if $X$ and $\psi_{0}$ are nontrivial), and the q-th power Frobenius should act on it by $g\left(\tilde{X}, \psi_{0}{ }^{\circ} T r\right)$, where


Since $G$ acts on $A^{+}(U)=R \ll x, y \gg /\left(y^{P}-y-x^{d}\right)$, we can decompose

$$
A^{\dagger}(J) \otimes K=\stackrel{\text { all }}{ }_{\oplus}^{X}, \psi_{0}\left(A^{\dagger}(U) \otimes K\right)^{X \times \psi_{0}}
$$

The following proposition is due to Monsky. We let $X_{a}$ denote $\xi \mapsto \omega(\xi)^{\frac{a}{d}(q-1)}$ $\qquad$ $\psi_{\text {triv }}$ denote the trivial character on $F_{p}$, and we let $\psi_{\pi}$ denote the character $a \longmapsto E_{\pi}(\omega(a))$ on $F_{p}$ (see the beginning of $\S 5$ ).

Proposition. The subspace of $A^{\dagger}(U)$ invariant under
$\{1\} \times Z / \mathrm{pZ} \subset \in$ is $\mathrm{R}\langle<\mathrm{x}\rangle>$; the subspace of $A^{\dagger}(\mathrm{U}) \otimes \mathrm{K}$ on which $2 / \mathrm{PZ}$ acts by $\psi_{\pi}$ is $E_{\pi}(y)(R<\langle x\rangle>\otimes R) ; \quad A^{+}(U)^{X_{a}} \times \psi_{\text {triv }}=$
$=x^{a} R \ll x^{d} \gg ;$ and $\left(A^{+}(U) \otimes K\right)^{X_{a}^{x} \psi_{\pi}}=E_{\pi}(y) x^{a}\left(R \ll x^{d} \gg \otimes K\right)$.
To prove this, first note that $E_{\pi}(y)$ and $E_{\pi}(-y)=1 / E_{\pi}(y)$ have coefficients in $R$ (see SI.3) and converge on a disc strictly larger than $D(1)$; hence $E_{\pi}(y)$ is a unit in $R<\langle y \gg$. In additiori, $E_{\pi}(y)$ transforms by $\psi_{\pi}$ under the action of $2 / \mathrm{p} Z$, i.e.,

$$
\begin{equation*}
\phi_{\alpha}\left(E_{\pi}(y)\right)=\psi_{\pi}(\alpha) E_{\pi}(y), \quad \alpha \in F_{p} \tag{6.1}
\end{equation*}
$$

where $\phi_{\alpha}$ is the Iifting of $y \longmapsto y+\alpha$ constructed in §5. To see this, note that $\phi_{\alpha}\left(E_{\pi}(y)\right) / E_{\pi}(y)$ is a $p-t h r$ root of unity (independent of $y$ ), because

$$
\left(E_{\pi}\left(\phi_{\alpha}(y)\right) / E_{\pi}(y)\right)^{p}=e^{p \pi\left(\phi_{\alpha}(y)-\phi_{\alpha}(y)^{p}\right)} / e^{p \pi\left(y-y^{p}\right)}=1
$$ since $\phi_{\alpha}(y)^{p}-\phi_{\alpha}(y)=y^{p}-y$. (Note that the extra $p$ in the exponent allows us to evaluate the exponent first, as at the beginning of §5.) To determine which p-th root, we set $y=0$ :

$$
E_{\pi}\left(\phi_{\alpha}(0)\right) / E_{\pi}(0)=E_{\pi}\left(\phi_{\alpha}(0)\right)=E_{\pi}(\omega(\alpha))=\psi_{\pi}(\alpha)
$$

and (6.1) is proved.
Since $A^{+}(J) \otimes R$ clearly has rank $p$ over $R \ll x \gg \otimes R$ (it
equals $\left.\underset{0 \leq i<p}{\oplus} y^{i} R \ll x \gg(0) K\right)$, and the subspace on which $z / p Z$ acts by $\psi_{\pi}$ (resp. $\psi_{\text {triv }}$ ) inciudes $\left.E_{\pi}(y) R<\langle x\rangle\right\rangle \otimes \mathbb{R}$ (resp. $R \ll x \gg(\times)$ ), the first two assertions of the proposition follow. The assertions about the action of $X_{a}$ are obvious.

Remarks. 1. For fixed $\pi, \pi^{p-1}=-p$, we could also use the units $E_{\pi}(y)^{i} \in R<\left\langle y \gg, i=0,1, \ldots, p-1\right.$, since $E_{\pi}(y)^{i}$ transforms by $\psi_{\pi}{ }^{\ddagger}$ under $\mathrm{z} / \mathrm{pZ}$. Hence,
$A^{+}(U) \otimes K=\underset{0 \leq i<p-1}{\oplus} E_{\pi}(y)^{i}(R \ll x \gg(\otimes K)$.
Note that $E_{\pi}(y)^{p}=e^{p \pi\left(y-y^{p}\right)}=e^{-p \pi x^{d}} \in R\langle\langle x\rangle\rangle$. Thus, $E_{\pi}(y)$ is a Kummer generator of $A^{\dagger}(U) \otimes K$ over $R \ll x \gg \otimes K$.
2. All of this applies to any curve of the form $y^{p}-y=f(x)$ $\in R[x]$; the specific polynomial $f(x)=x^{d}$ was not needed in analyzing the action of $\mathrm{z} / \mathrm{pz}$.

Proposition. For $\psi_{0}=\psi_{\pi}$ and for $X=\chi_{a}, 1 \leq a<d$,
$H^{1}\left(A^{\dagger}(\mathrm{U}) \otimes \mathrm{K}\right) \mathrm{X}^{\times} \psi_{0}$ is a one-dimensional $K$-vector space with basis $E_{\pi}(y) x^{a} \frac{d x}{x} ; \quad$ if either $X$ or $\psi_{0}$ is trivial, this space is zero.

Proof. We compute: $d E_{\pi}(y)=E_{\pi}(y) d \log E_{\pi}(y)=$ $E_{\pi}(y) d\left(\pi\left(y-y^{p}\right)\right)=\pi E_{\pi}(y) d\left(-x^{d}\right)=-\pi d x^{d-1} E_{\pi}(y) d x$. Thus,
$d\left(E_{\pi}(y) x^{a} f\left(x^{d}\right)\right)=E_{\pi}(y) x^{a d x} \frac{x}{x}\left(-\pi d x^{d} f\left(x^{d}\right)+a f\left(x^{d}\right)+d x^{d} f^{\prime}\left(x^{d}\right)\right)$, so that for $a \geq 1$,

$$
\begin{aligned}
\left(H^{1}\left(A^{+}(\mathrm{U}) \otimes \mathrm{K}\right)\right)^{\chi \times \psi_{0}} & =\frac{E_{\pi}(y) x^{a} R \ll x^{d}>\frac{d x}{x} \otimes K}{d\left(E_{\pi}(y) x^{a} R \ll x^{d} \gg \otimes K\right)} . \\
& \simeq \frac{R \ll x \gg \otimes R}{\left(-\pi x+\frac{a}{d}+\frac{d}{d x}\right)(R \ll x \gg \otimes K)} .
\end{aligned}
$$

The proposition now asserts that the cokernel of $\partial=x \frac{d}{d x}+\frac{a}{d}-\pi x$
on $R \ll x \gg \otimes K$ is $K$ if $a \geq 1$. (The triviality of $\mathrm{H}^{1}\left(\mathrm{~A}^{+}(\mathrm{U}) \otimes \mathrm{K}\right) \mathrm{X}^{\times} \psi_{\text {triv }}$ is immediate, since $\frac{d}{\mathrm{dx}}$ is surjective on $R \ll x \gg \otimes K$; the triviality of $H^{1}\left(A^{\dagger}(U) \otimes K\right){ }^{X_{\text {triv }} \times \psi_{0}} \quad$ will be shown later.)
We compute:

$$
x^{m+1}=\frac{m+\frac{a}{d}}{\pi} x^{m}+\partial\left(-\frac{x^{m}}{\pi}\right)
$$

so that

$$
\begin{aligned}
\sum_{m=0}^{\infty} b_{m+1} x^{m+1}= & \sum_{m=0}^{\infty} b_{m+1} \frac{\left(m+\frac{a}{d}\right)\left(m+\frac{a}{d}-1\right) \cdots\left(\frac{a}{d}\right)}{\pi^{m+1}} \\
& -\partial\left(\sum_{n=1}^{\infty} x^{n} \sum_{m \geq n} b_{m+1} \frac{\left(m+\frac{a}{d}\right)\left(m+\frac{a}{d}-1\right) \cdots\left(n+1+\frac{a}{d}\right)}{\pi^{m+1-n}}\right)
\end{aligned}
$$

## Because we have

$$
\begin{aligned}
\operatorname{ord}_{p} b_{m+1} \frac{(m+1)!}{\pi^{m+1}}\binom{m+\frac{a}{d}}{m+1} & \geq \text { ord }_{p} b_{m+1} \frac{(m+1)!}{\pi^{m+1}} \quad \begin{array}{l}
\text { (sees } \delta I .3, \\
\text { formula (3.3)) }
\end{array} \\
& \left.=\text { ord }_{p} b_{m+1}-\frac{S_{m+1}}{p-1} \quad \text { (see } \S I .3,(3.1)\right),
\end{aligned}
$$

while ord $b_{m+1} \geq(m+1) \varepsilon+B$ when $\left.\sum b_{m+1} x^{m+1} \in R \ll x\right\rangle>\otimes R$, it follows that the constant sum in (6.2) converges (since $S_{m+1} \leq$ ( $\left.\mathrm{p}-1)\left(\log _{\mathrm{p}} \mathrm{m}+1\right)\right)$. Similarly, ord ${ }_{p}$ of the inner sum inside the
$\partial$ is at least

$$
\min _{m \geq n}\left((m+1) \varepsilon+B-\frac{1+S_{m-n}}{p-1}\right) \geq n \varepsilon^{r}+B^{\prime}
$$

for some $\varepsilon^{\prime}>0$ and some real number $B^{\prime}$. Thus, any element $f$
in $R<\langle X\rangle>\otimes \mathbb{K}$ can be written as const $+\partial g, g \in R<\langle x\rangle>\otimes K$. It remains to show that 1 cannot be written as $\partial g, g \in R \ll x \gg \otimes R$, if $a \geq 1$, in which case we will have proved that the cokernel is precisely the constants. Suppose $g=\sum \mathrm{b}_{\mathrm{n}} \mathrm{x}^{\mathrm{n}}, \quad 1=\partial \mathrm{g}$. Then

$$
I+\pi x \sum b_{n} x^{n}=\left(x \frac{d}{d x}+\frac{a}{d}\right) \sum b_{n} x^{n},
$$

and comparing coefficients gives
$1=\frac{a}{d} b_{0}$

$$
\begin{aligned}
b_{n} & =\frac{\pi}{n+\frac{a}{d}} b_{n-1} \quad(n \geq 1) \\
& =\cdots=\frac{\pi^{n}}{(n+1)!\binom{n+\frac{a}{d}}{n+1}} .
\end{aligned}
$$

Thus, $\operatorname{ord}_{p} b_{n} \leq\left(S_{n+1}-1\right) /(p-1)$, which is not $\geq \varepsilon_{n}+B$ for any $\varepsilon>0$ and real number B. Thus, $g \notin R<\langle x\rangle>(\otimes) K$.

The final assertion is that the cohomology is zero when $a=0$. Roughly speaking, this is because the basis $E_{\pi}(y) x^{a} \frac{d x}{x}$ for the cokernel of $\partial$ only makes sense when $a \geq 1$. More precisely,

$$
\begin{aligned}
& H^{1}\left(A^{+}(J) \otimes R\right) X_{\text {triv }} \times \psi_{\pi}=\frac{E_{\pi}(y) x^{d-1} \cdot R^{<} \ll x^{d} \gg d x \otimes K}{d\left(E_{\pi}(y) R \ll x^{d} \gg(\otimes K)\right.} \\
& \simeq \frac{R<\langle x\rangle>\otimes) K}{\left(\frac{d}{d x}-\pi\right)(R \ll x \gg \otimes \otimes)} \text {. }
\end{aligned}
$$

To show that $\frac{d}{d x}-\pi$ is surjective on $\mathbb{R}\langle\langle x\rangle\rangle \otimes \mathbb{R}$, it suffices to show that its formal inverse $-\frac{1}{\pi}\left(1-\frac{1}{\pi} \frac{d}{d x}\right)^{-1}=-\sum_{n=0}^{\infty} \pi^{-n-1}\left(\frac{d}{d x}\right)^{n}$ takes a series in $R<\langle x\rangle>\otimes K$ to a series in $R<\langle x\rangle>\otimes \pi$. But this is easy using the same estimate for $\frac{n!}{\pi^{n}}\binom{\alpha}{n}$ as above. This completes the proof of the proposition.

Recall that our Gauss sum $g(\tilde{X}, \psi), \psi=\psi_{0}{ }^{\circ} T r$, is the action of the q -th power Frobenius $F=F_{\mathrm{F}_{\mathrm{q}}}$ on $\mathrm{H}^{1}\left(\mathrm{~A}^{+}(\mathrm{U}) \otimes \mathrm{K}\right){ }^{\mathrm{X}}{ }^{\times} \psi_{0}$. Since the curve $A(d, p)$ is defined over $F_{p}$, the $p$-th power Frobenius $F_{0}=F_{\mathrm{F}}$ also acts on $\mathrm{A}(\mathrm{d}, \mathrm{p})$ and hence on $\mathrm{H}^{1}\left(\mathrm{~A}^{+}(\mathrm{U}) \otimes \mathrm{K}\right)$; and we have $F=F_{0}^{f}$, where $q=\mathrm{p}^{\mathrm{f}} \equiv 1$ (mod d). However, $F_{0}$ does not commute with the action of $G=\mu_{d} \times Z / \mathrm{p} Z$; its matrix in the eigen vectors of $G$ is not diagonal. More precisely, we have

Proposition. $\quad F_{0}^{*}\left(\mathrm{H}^{1}\left(\mathrm{~A}^{+}(\mathrm{U}) \otimes \mathrm{K}\right)^{X \times \psi_{0}}\right) \subset \mathrm{H}^{1}\left(\mathrm{~A}^{+}(\mathrm{J}) \otimes \mathrm{K}\right) \mathrm{X}^{\mathrm{p} \times \psi_{0}} 0$.
This is an immediate consequence of the commatation relation $F_{0} \circ(\xi, \alpha)(\mathrm{x}, \mathrm{y})=\left(\xi^{\mathrm{P}} \mathrm{x}^{\mathrm{p}}, \mathrm{y}^{\mathrm{p}}+\alpha^{\mathrm{P}}\right)=\left(\xi^{\mathrm{P}} \mathrm{x}, \mathrm{y}^{\mathrm{P}}+\alpha\right)=\left(\xi^{\mathrm{P}}, \alpha\right) \circ F_{0}(\mathrm{x}, \mathrm{y})$, for $(\xi, \infty) \in G$.

$$
\text { Thus, if } X=X_{a}, l \leq a<d \text {, and if we let } a^{\prime}, a^{\prime \prime}, \ldots, a^{(j)} \text {, }
$$

$\ldots, a^{(f-1)}$ denote the least positive residue of $p^{j} a \bmod d$,
then for some constant $\lambda=\lambda(a, d, \pi)$ we have
$F_{\hat{0}}^{*}\left(E_{\pi}(y) x^{a} \frac{d x}{x}\right)=\lambda E_{\pi}(y) x^{a^{\prime}} \frac{d x}{x}$ in $H^{1}\left(A^{\dagger}(U) \otimes K\right) . \quad$ (6.3)
Then we obtain

$$
\begin{equation*}
g\left(\tilde{X}_{a}, \psi_{\pi} \pi^{\circ} T r\right)=\prod_{j=0}^{\mathrm{f}-1} \lambda\left(a^{(j)}, \mathrm{d}, \pi\right) . \tag{6.4}
\end{equation*}
$$

Proposition. $\lambda(a, d, \pi)=-p \pi^{-\frac{p a-a^{i}}{d}} \Gamma_{p}\left(1-\frac{a^{d}}{d}\right)$.
This proposition, together with (6.4) and the elementary fact that $\frac{p-1}{d} \sum_{j=0}^{f-1} a^{(j)}=$ the sum of the $p$-adic digits in $a \frac{q-1}{d}$, imply:

Theorem.

$$
\left.\left.g\left(\tilde{x}_{a}, \psi_{\pi} \cdot \mathrm{Tr}\right)=\frac{(-\mathrm{p})^{\ddagger} \prod_{j=0}^{f-1} \Gamma_{p}\left(1-\left\langle\frac{p^{j}}{d}\right\rangle\right.}{\left.\mathrm{S}_{\mathrm{a}}\right\rangle}\right\rangle\right),
$$

where $<>$ denotes the least positive residue modulo 1 (not to be confused with our earlier meanings of < $>$ ).

$$
\begin{align*}
& \text { Proof of proposition. Let } R \text { be the ring of integers in } \\
& Q_{p}(\sqrt{d} \sqrt{1}, \pi) \text {. First we imbed } \\
& A^{+}(\mathbb{U})=R \ll x, y \gg /\left(y^{p}-y-x^{d}\right) \xrightarrow{\phi} R[[x]] \tag{6.5}
\end{align*}
$$

by sending $y$ to the formal power series solution of $y^{p}-y=x^{d}$ near $(0,0)$, i.e., with zero constant term: $y=-x^{d}+\ldots$. To see that such a solution exists (and is unique), one can use a version of Hensel's lemma (5I.2) for the "x-adic topology" in R[[x]] (the topology which says that two series are close together if their difference is divisible by a large power of $x$ ). Then $y=0$ is a solution $\bmod x$ of $f(y)=y^{p}-y-x^{d}=0$; moreover $f^{\prime}(0)=$ $-1 \neq 0$, so the existence and uniqueness of the desired series $\phi_{y}(x) \in R[[x]]$ is assured.

The imbedding (6.5) induces a homomorphism
$\left.\mathrm{H}^{1}\left(\mathrm{~A}^{\dagger}(\mathrm{U}) \otimes \mathrm{K}\right) \longrightarrow \mathrm{H}^{1}(\mathrm{R}[\mathrm{X}]] \otimes \mathrm{K}\right)$
which commutes with the p -th power Frobenius $F_{0}^{*}$, where $F_{0}$ acts on $R[[x]]$ by sending $x$ to $x^{p}$ and acts on $A^{\dagger}(U)$ by sending $x$ to $x^{p}$ and $y$ to the unique element which lifts $y^{p}$ and satisfies $F_{0}(y)^{\mathrm{p}}-F_{0}(\mathrm{y})=\mathrm{x}^{\mathrm{pd}}$.

## We have

$$
\phi\left(E_{\pi}(y) x^{a} \frac{d x}{x}\right)=\phi\left(e^{\pi\left(y-y^{p}\right)} x^{a} \frac{d x}{x}\right)=e^{-\pi x^{d}} x^{a} \frac{d x}{x}
$$

Thus, by (6.3),

$$
\begin{equation*}
F_{0}^{*}\left(e^{-\pi x^{d}} x^{a} \frac{d x}{x}\right)=\lambda(a, d, \pi) e^{-\pi x^{d}} x^{a^{\prime}} \frac{d x}{x} \tag{6.6}
\end{equation*}
$$

in $H^{1}(R[[x]] \otimes R)$. In other words, the difference between the two sides of this equality lie in $d(R[f x]] \otimes R)$. Note that
$\sum_{n=1}^{\infty} a_{n} x^{n} \frac{d x}{x}$ is in $d(R[[x]]) \Leftrightarrow \frac{a_{n}}{n} \in R$ for all $n$.
Thus,

$$
\sum_{n=1}^{\infty} a_{n} x^{n} \frac{d x}{x} \in d(R[[x]] \otimes R) \quad \Leftrightarrow \quad \operatorname{ord}_{p} a_{n} \geq \operatorname{ord}_{p}{ }^{n}+\text { constant }
$$

$$
\text { Evaluating } F_{0}^{*} \text { in (6.6) gives }
$$

$$
\begin{equation*}
e^{-\pi x^{p d}} x^{p a} p \frac{d x}{x} \equiv \lambda e^{-\pi x^{d}} x^{a^{\prime}} \frac{d x}{x} \quad \bmod d(R[[x]] \otimes R) \tag{6.7}
\end{equation*}
$$

Equating powers of $x$ in (6.7) gives

$$
\begin{equation*}
\frac{(-\pi)^{n}}{n!} p \equiv \lambda \frac{(-\pi)^{m}}{m!} \quad \bmod \quad p^{\text {ord }_{p}\left(m d+a^{\prime}\right)+\text { const }}, \tag{6.8}
\end{equation*}
$$

where $m$ is chosen so that $m d+a^{r}=p(n d+a)$, i.e., $m=$ $\frac{p a-a^{1}}{d}+p n$.

We now choose a sequence of $n^{\prime} s$ for which ord $(n d+a)$ (and hence $\operatorname{ord}_{p}\left(m d+a^{\prime}\right)$ ) approaches infinity. Namely, we can let $n$
approach $-\mathrm{a} / \mathrm{d}$ by taking the p -adic expansion of $-\mathrm{a} / \mathrm{d}$; but it is more convenient to expand $-a / d$ in powers of $q=p^{f}$. Letting $b=\frac{q-1}{d} a$, so that $0<b<q-1$, we have

$$
-\frac{a}{d}=\frac{b}{1-q}=b+b q+b q^{2}+\ldots
$$

Let $n=b+b q+\ldots+b q^{j-1}=-\frac{1-q^{j}}{d} a$. Then ord $(n d+a)$ $\geq$ fj. Further note that ord $\left(\frac{n!}{\pi^{n}}\right)=-\frac{s_{n}}{p-1}=-\frac{j S_{b}}{p-1}$. Similarly, $\operatorname{ord}_{p}\left(\frac{m!}{\pi^{m}}\right):=-\frac{j S_{b}}{p-1}+$ const. Thus, multiplying (6.8) through by $m!(-\pi)^{-m}$ and carefully taking note of ord $p$ and the effect on the congruence, we obtair

$$
\lambda \equiv p(-\pi)^{n-m} \frac{m!}{n!} \bmod p j-j S_{b} /(p-1)+\text { const } .
$$

But $\mathrm{j}\left(\mathrm{f}-\frac{\mathrm{S}_{\mathrm{b}}}{\mathrm{p}-1}\right) \longrightarrow \infty$ as $\mathrm{j} \longrightarrow \infty$ (this is because $\mathrm{b}<\mathrm{q}-1$ and so has at least one digit less than $p-1$; thus, $S_{b}<f(p-1)$ ). Hence,

$$
\lambda=p \lim _{j \rightarrow \infty}(-\pi)^{n-m} \frac{m!}{n!}
$$

where $n=\frac{a}{d}\left(q^{j}-1\right), \quad m=p n+\frac{p a-a^{\prime}}{d}$. Note that $\quad \prod_{i \leq m, p \mid i} i=$ $p \cdot 2 p^{\circ} \cdot n p=p^{n} n!$.

Thus, by the definition of $\Gamma_{p}$,

$$
\Gamma_{p}(m+1)=(-1)^{m+1} \prod_{i \leq m, ~ p q_{i}} i=(-1)^{m+1} \frac{m!}{n!p^{n}}
$$

Hence,

$$
\lambda=p \lim _{j \rightarrow \infty}(-\pi)^{n-m}(-1)^{m+1} p^{n} \Gamma_{p}(m+1)
$$

Now

$$
(-\pi)^{n-m}(-1)^{m+1} p^{n}=-\pi^{n-m}(-p)^{n}=-\pi^{p n-m},
$$

since $\pi^{p-1}=-p$. Since $p n-m=-\frac{p a-a '}{d}$; and $n \longrightarrow-\frac{a}{d}$ and $m \longrightarrow-\frac{a^{j}}{d}$ as $j \longrightarrow \infty$, we conclude:

$$
\lambda=-p \pi^{-\left(p a-a^{\prime}\right) / d} \lim _{m \rightarrow-a^{\prime} / d} \Gamma_{p}(m+1),
$$

and the proposition is proved.
Remarks. 1. For simplicity, suppose $p \equiv I(\bmod d)$. Then the theorem reads:

$$
\begin{equation*}
g\left(\tilde{x}_{a}, \psi_{\pi}\right)=-p \pi^{-a(p-1) / d} \Gamma_{p}\left(1-\frac{a}{d}\right) . \tag{6.9}
\end{equation*}
$$

Suppose $1 \leq r, s, r+s<d$, and let $\bar{X}_{a}$ denote $\tilde{X}_{d-a}$. By property (3) of Gauss and Jacobi sums at the beginning of this chapter, (6.9) gives us:

$$
J\left(\bar{\chi}_{r}, \bar{X}_{s}\right)=\frac{\Gamma_{p}\left(\frac{r}{d}\right) r_{p}\left(\frac{s}{d}\right)}{\Gamma_{p}\left(\frac{r+s}{d}\right)}
$$

which looks remarkably similar to the beta function value in (2,2) for the classical periods of the differential $\omega_{r, s}$ !
2. The above theorem also gives an analogy between the ChowlaSelberg formula for the periods of an elliptic curve with complex multiplication and a p-adic expression for the roots of the zeta function of the elliptic curve; see [37,34].

We have thereby shown Gauss sums to be p-adic analogs of special values of the gamma function. In the next section, we show how Stickleberger's theorem on the ideal decomposition of Gauss sums is an immediate corollary. In the next chapter we shall see a subtler application: the proof that the p-adic Dirichlet L-function $L_{p}(s, X)$ has at most a simple zero at $s=0$

## 7. Stickleberger's theorem

Stickleberger's theorem gives the ideal decomposition of Gauss sums $g(\chi, \psi)$ in $\bar{Q}$. Let $K$ denote the field $Q(\sqrt{d} \sqrt{1})$. Let $P$ be any fixed prime ideal of $K$ lying over $p$. Let $q$ be the number of elements in the residue field $0 / P$ of $P$; thus, $q=p^{f}$ is the least power of $p$ such that $d \mid q-1$. We identify $F_{q}$ with $0 / P$, and let

$$
X: \mathrm{F}_{\mathrm{q}}^{*}=(0 / \mathrm{P}) * \longrightarrow \mu_{\mathrm{d}} \subset \mathrm{~K}
$$

be a multiplicative character. Let $a$ be the integer, $0 \leq a<d$, determined by $\chi(x) \bmod P=x^{a(q-1) / d}$ for $x \in 0 / P$. Thus, if we use $P$ to imbed $K$ in $\Omega_{p}$, and consider $\chi$ to take values in $\Omega_{p}$, then $X$ is the $a \frac{q-1}{d}$ - th power of the Teichmuller character, i.e., $x=\tilde{X}_{a}$ in the notation of $\S 53-6$.

The Gauss sumn $g(X, \psi)$ is obviously an algebraic integer in $\mathrm{K}(\sqrt{\mathrm{p}} \sqrt{1})$. By checking the action of $\mathrm{GaI}(\mathrm{K}(\sqrt[\mathrm{P}]{1}) / \mathrm{K})$ on
$g(X, \psi)=-\sum_{x \in F_{q}^{*}} \chi(x) \psi(x)$,
we see that $g(\chi, \psi)^{\mathrm{d}}$ lies in K and is independent of the additive character $\psi$. By property (1) in $\S 1, g(\chi, \psi)$ divides $q=p^{f}$; hence, the ideal $\left(g(X, \psi)^{d}\right)$ in $K$ must be a product of powers of prime ideals of $K$ which divide $p$. Stickleberger's theorem gives these powers.

We can consider $g(X, \psi)^{d} \quad$ p-adically if we choose an imbedding $\tau: K \longrightarrow Q_{p}(\sqrt[d]{1})$. As explained before ( $\$ 11,2$ ), for our fixed prime ideal $P$ of $K$ dividing $p$, we obtain such an imbedding $l^{=} i_{p}$ by taking the completion of $K$ in the P-adic topology; this imbedding $I$ allows us to identify $X$ (strictly speaking,
$20 \chi$ ) with $\omega^{\frac{a}{d}(q-1)}$. The power of $P$ dividing the ideal $\left(g(\chi, \psi)^{d}\right)$ in $K$ is simply ord $\left(i_{p}\left(g(\chi, \psi){ }^{d}\right)\right)$.

Recall that $\operatorname{Gal}(K / Q) \simeq(Z / d Z) *$, where $\sigma_{j} \vdots \xi \longrightarrow \xi^{j}$ for $j \epsilon(Z / d Z)^{*}, \quad \xi^{d}=1$. Gal(K/Q) permutes the prime ideals of $K$ dividing $P$, and we let $P_{j}=\sigma_{j} P$. of course, $P_{j}=P_{j}$, if $j / j$ ' is in the subgroup of powers of $p$ in ( $Z / \mathrm{dZ}$ )* (the "decomposition group" of P).

If $a / b \in Q$ with g.c.d. $(b, d)=1$, let $\langle a / b\rangle{ }_{d}$ denote the least positive residue of $a / b$ mod $d$, i.e., the least positive
k such that $\mathrm{kb} \equiv \mathrm{a}(\bmod \mathrm{d})$. Let

$$
\left(g(x, \psi)^{d}\right)=\prod_{j \in(z / d z) * /\left\{p^{k}\right\}_{j}} p_{j}^{\alpha}
$$

be the ideal decomposition of $g(\chi, \psi)^{d}$.

$$
\text { Stickleberger's theorem. } \quad \alpha_{j}=\sum_{k=0}^{\mathrm{f}-1}\left\langle-a / j p^{k^{\prime}}\right\rangle_{d} \text {, i.e., }
$$

$$
\left(g(x, \psi)^{d}\right)=p_{\sigma}\left(\sum_{j \in(z / d Z) *}\langle-a / j\rangle{ }_{d} \sigma_{j}\right)^{k=},
$$

where we write $P^{\sigma_{j}}$ for $\sigma_{j} P=P_{j}$.

$$
\text { Proof. } \alpha_{j}=\text { power of } P_{j} \text { dividing } g(\chi, \psi)^{d}
$$

$=$ power of $P$ dividing $\sigma_{j}^{-I} g(X, \psi)^{d}$
$=$ power of $P$ dividing $g\left(\chi^{\left(j^{-1}\right)}, \psi\right)^{d}$,
where $X^{\left(j^{-1}\right)}=\sigma_{j}^{-1} \circ \chi$ is the character $X^{\langle 1 / j\rangle} d=\omega^{\langle a / j\rangle} d \frac{q-1}{d}$. But according to the theorem in $\S 6$,
$\operatorname{ord}_{p}\left(i_{p}\left(g\left(\chi^{\left(j^{-1}\right)}, \psi\right)^{d}\right)\right)=\sum_{k=0}^{f-1}\left(d-\left\langle p^{k} a / j\right\rangle_{d}\right)=\sum_{j=0}^{f-1}\left\langle-a / j p^{k}\right\rangle_{d} . \quad$ Q.E.D.

## 1. Regulators and L-functions

If $K$ is a number field with $r_{1}$ real imbeddings and $2 r_{2}$ complex imbeddings, a classical theorem of pirichlet [13] asserts that the multiplicative group $E$ of units of $K$ is the direct product of the (finite) group $W$ of roots of 1 in $K$ and a free abelian group of rank $r_{1}+r_{2}-1$, i.e., there exist units $e_{1}, \ldots$, $e_{r_{1}+r_{2}-1}$ such that every unit can be written uniquely in the form

$$
\sum_{1}^{T I} e_{1}^{-1} e_{2}^{m_{2}} \cdots e_{r_{1}+r_{2}}^{m_{r_{1}}+x_{2}-1}, \quad m_{i} \in Z, \quad \eta \text { a root of } 1
$$

If $\phi_{1}, \ldots, \phi_{r_{1}}$ denote the real imbeddings and $\phi_{r_{1}+1}, \ldots, \phi_{r_{1}}+r_{2}$ denote the complex imbeddings (one chosen from each complex conjugate pair), then the map

$$
\begin{aligned}
& \left(\log \left|\phi_{1}()\right|, \ldots, \log \left|\phi_{r_{1}}()\right|,\right. \\
& \left.\quad 2 \log \left|\phi_{r_{1}+1}()\right|, \ldots, 2 \log \left|\phi_{r_{1}+r_{2}}()\right|\right): K \longrightarrow R^{r_{1}+r_{2}}
\end{aligned}
$$

takes the group $E / W$ of units modulo roots of 1 isomorphically to a lattice in the hyperplane $x_{1}+x_{2}+\cdots+x_{r_{1}}+r_{2}=0$ in $r_{1}+r_{2}$ dimensional real space. The volume

$$
\operatorname{det}\left(n_{i} \log \left|\phi_{i}\left(e_{j}\right)\right|\right), \quad 1 \leq i, j \leq r_{1}+r_{2}^{-1}
$$

where $n_{i}=1$ for $i \leq r_{1}, n_{i}=2$ for $i>r_{1}$, of a fundamental parallelotope (actually of its projection on the $x_{r_{1}+r_{2}}{ }^{\text {-hyperplane) }}$ is called the (classical) regulator of K . It depends only on K ,
not on the choice of "fundamental units" $e_{j}$ or the ordering of the $\phi_{i}$, and it is always nonzero.

In this chapter we discuss two very different types of p-adic regulators. The first type, due to Leopoldt, takes the same units $e_{j}$ as in the classical case, replaces the imbeddings $\phi_{i}: K \hookrightarrow C$ by imbeddings $\phi_{i}$ of K into the algebraically closed field $\Omega_{p}$, and replaces $\log$ by $\ln _{p}$. Because there is no natural way of eliminating $r_{2}$ of the p-adic $\phi_{i}$ the way we eliminated one $\phi_{i}$ from each complex conjugate pair in the complex case, Leopoldt further assumes that $r_{2}=0$, i.e., $K$ is totally real.

The second, more recently developed type of p-adic regulator is due to B. H. Gross ([35], see also Greenberg [31]). It applies to number fields K which are totally complex, i.e., $\mathrm{r}_{1}=0$; more specifically, to so-called CM fields, which are quadratic imaginary extensions of a totally real field $K^{+}$, i.e., $K=K^{+}(\sqrt{-\alpha})$ for some $\alpha$ which is positive under all imbeddings $\mathrm{K}^{+} \hookrightarrow \mathrm{R}$. Gross works not with ordinary units, but rather with "p-units". A p-unit is an algebraic number which has absolute value 1 at all places (including archimedean valuations) except for those over. p. In other words, all of its conjugates must have complex absolute value 1 , and its ideal decomposition only involves primes dividing $p$. A key example of a p-unit is $g / \bar{g}$, where $g$ is a Gauss sum for a finite field of $q=p^{f}$ elements. The multiplicative group of p-units, if tensored with $Q$ and written additively, is isomorphic to the vector space of divisors $\sum_{\left.p\right|_{P}} Q(P-\bar{P})$. Anong the imbeddings $\phi_{i}: K \longrightarrow \Omega_{p}$, Gross chooses one from each coset modulo $\pm$ (the decomposition group of $p$ ), so that each $\phi_{i}$ gives a different permutation of the divisors $P-\bar{P}$. For a more precise statement, see below. Gross then takes the determinant of a $\frac{g}{2} \times \frac{g}{2}$ matrix, where $g$ is the number of primes $P$ over $p$. Gross's padic regulator is very different from Leopoldt's. In Gross's case the set of units considered and even the size of the matrix vary completely from one $p$ to another.

The basic way in which regulators occur "in nature" is in the expansion at 1 or 0 of $\zeta$ - and L-functions. First, in the classical case, let

$$
\zeta_{K}(s)=\sum \frac{1}{(N A)^{s}}
$$

be the Dedekind zeta function of the number field $K$. Here the sum is over all non-zero integral ideals $A$ of $K$, and $N$ is the norm. The series converges for Res $>1$ and can be analytically continued to a function which is holomorphic on the complex plane except for a simple pole at $s=1$. The residue at $s=1$ equals

$$
\frac{2^{r_{1}}(2 \pi)^{r_{2} h R}}{w \sqrt{|D|}},
$$

where $r_{1}$ and $r_{2}$ are, as above, the number of real imbeddings and pairs of complex imbeddings; $w$ is the number of roots of 1 in $K, D$ is the discriminant of $K$, $h$ is its class number, and R is its regulator. The subtlest and most elusive term in this formula is the regulator.

The Leopoldt p-adic regulator occurs in a similar way. Let $K$ be a totally real field, i.e., $r_{2}=0$. Serre [85] has shown how to associate to K a p-adic zeta function $\zeta_{\mathrm{K}, \mathrm{p}}(\mathrm{s})$ which is defined and holomorphic on the closed unit disc in $\Omega_{p}$ (actually, on a slightly larger disc) except for a possible pole at $s=1$. Conjecturally, the pole at $s=1$ is a simple pole with residue given by

$$
\frac{2^{r_{1}} h_{p, \text { Leopoldt }}}{w \sqrt{D}} \text {. }
$$

Here all of the terms have been defined above except for $E$, which is a product of Euler factors. (The phenomenon of "throwing out the p-Euler factor" can be expected to occur in all p-adic versions cf classical formulas for $\zeta$ - and L-functions.) In the simplest case, when $K=Q$, we have $r_{1}=1, h=1, w=2, D=1, R_{p, \text { Leopoldt }}$ $=1$, and $E=1-\frac{1}{p}$.

Note that there's an ambiguity of sign in $\sqrt{\mathrm{D}}$. We will see that also $R_{p}$, Leopoldt is only defined up to a sign. In the
classical case one can normalize by taking the absolute value of the determinant in the definition of $R$ and the positive square root of $|\mathrm{D}|$. It is harder to $\mathrm{f} x \mathrm{x}$ the sign in the p-adic case.

It is also conjectured that always $R_{p, \text { Leopoldt }} \neq 0$, i.e., there really is a pole at $s=1$. Both the residue formula and the non-vanishing of the regulator have been proved in the case when $K$ is an abelian extension of $Q$ (the "abelian over $Q$ " case). In that case $\zeta_{\mathrm{K}}$ is a product of Dirichlet L-series, and the necessary facts were essentiaily worked out by Leopoldt [64] (see also [61]). We shall prove the non-vanishing of $R_{p, \text { Leopoldt }}$ below in the abelian over $Q$ case. The proof uses a theorem from transcendence theory, which will be stated without proof.

But very little is known about $R_{p}$, Leopoldt in the non-abelian case. A partial result supporting the residue formula was obtained by Serre [86], who proved that for any totally real field $K$, if $\zeta_{\mathrm{K}, \mathrm{p}}$ has a pole at I , then $\mathrm{R}_{\mathrm{p}, \text { Leopoldt }} \neq 0$.

It should also be mentioned that the "Leopoldt conjecture" (non-vanishing of $R_{p}$, Leopoldt) and the expected relationship between the p-adic regulator and the residue at 1 has been generalized by Serre to p-adic "Artin L-functions" associated to representations of the Galois group of $\bar{k} / k$ ( $k$ totally real).

Gross's p-adic regulator, we shall see, is connected to the behavior near $s=0$ of $p$-adic Artin L-functions. These are p-adic L-functions $L_{p}(s, \rho)$ which p-adically interpolate values of the Artin L-series associated with a representation $\rho$ of the Galois group of a CM field $K$ over a totally real ground field $k$. The order of zero $m_{p}$ of $L_{p}(s, p)$ at $s=0$ has been conjectured for some time (see [29]). Gross further conjectures that the leading term in the Taylor series at 0 of $L_{p}(s, \rho)$ is
$R_{p, G r o s s}(\rho) A(\rho) \mathrm{s}^{m_{\rho}}$,
where $R_{p, G r o s s}(\rho)$ is Gross's $p$-adic regulator and $A(\rho)$ is an explicitly given algebraic number, which turns out to be a product
of certain Euler type factors and an algebraic number which is independent of p . (For a more detailed account, see below.) Note the analogy with the Leopoldt residue formula discussed above, in which the leading coefficient of the Laurent expansion (at $s=1$ ) is the product of a (p-adic transcendental) regulator, an Euler term, and an algebraic number independent of $p$. In the classical case, as we shall see below, the functional equation for L-functions gives a direct relationship between the expansion at $s=1$ and the expansion at $s=0$. But in the p-adic case there is no functional equation, and no one has yet been able to explain the analogy between the Leopoldt and the Gross formulas, in the sense of providing a direct link between the two types of p-adic regulators.

Gross's conjectured formula was motivated by: Ferrero-Greenberg's proof [29] that p-adic Dirichlet l-series have at most a simple zero at 0 ; and a conjecture of Stark and Tate concerning the leading coefficient at 0 of classical Artin L-series. Gross's conjecture is known to be true when $K$ is an abelian extension of $Q$. In the abelian over $Q$ case it reduces to the case when $\rho$ is a onedimensional character, $m_{\rho}=1$, and the conjecture asserts that
$L_{p}^{\prime}(s, \rho)=R_{p, G r o s s}(\rho) A(\rho) \neq 0$ at $s=0$.
We shall give Gross's variant of Ferrero-Greenberg's original proof of this fact.

Gross developed his conjecture as a p-adic analog of a conjecture of Stark [90] and Tate [93]. Instead of giving the Stark-Tate conjecture in the general setting, I'11 illustrate the idea by showing how it interprets the classical formula (see §II.5)

$$
\begin{equation*}
L(1, \chi)=-\frac{g_{X}}{d} \sum_{0<a<d} \bar{\chi}(a) \log \left(1-\zeta^{-a}\right) \tag{1.1}
\end{equation*}
$$

where $X$ is a nontrivial Dirichlet character of conductor $d, \zeta$ is a primitive d-th root of 1 , and $g_{\chi}=\Sigma \chi(a) \zeta^{a}$.

For simplicity, we take the case when $X$ is a nontrivial even character and $d=p^{N}$ is a power of an odd prime $p$.

Let $K=Q(\zeta)$, where $\zeta$ is a primitive d-th root of 1 , and let $K^{+}=Q\left(\zeta+\zeta^{-1}\right)$ be its maximal totally real subfield. Then $\operatorname{Gal}(\mathrm{K} / Q) \approx(Z / \mathrm{dZ}) *, \quad \sigma_{a}(\zeta)=\zeta^{a}$, and $\operatorname{Gal}\left(\mathrm{K}^{+} / Q\right) \approx(Z / \mathrm{dZ}) * /\{ \pm 1\}$. Let $G$ denote $(Z / d Z) * /\{ \pm 1\}$, so that summation over $a \in G$ means


Let $g$ be a generator of the cyclic group ( $Z / \mathrm{dZ}$ )* (recall that $d$ is an odd prime power), so that $\sigma_{g} \in \operatorname{Gal}\left(\mathrm{~K}^{+} / \mathrm{Q}\right)$ generates Ga1 ( $\left.\mathrm{K}^{+} / \mathrm{Q}\right)$. Let

$$
\varepsilon=\left(\zeta-\zeta^{-1}\right)^{\sigma}{ }^{-1}=\frac{\zeta^{g}-\zeta^{-g}}{\zeta-\zeta^{-1}}
$$

Then it can be shown ([61], p. 85) that $\varepsilon$ is a Minkowski unit in $K^{+}$(also in K), i.e., $\left\{\sigma_{a} \varepsilon\right\}_{a \epsilon G}$ generate a subgroup of finite index in the group $E$ of units of $K^{+}$. Let $E_{\text {cyc }} \subset E$ denote this group of "cyclotomic units". (Equivalently, the $\sigma_{a} \varepsilon$ are multiplicatively independent except for the single relation $\Pi_{a \in G} \sigma_{a} \varepsilon=N \varepsilon$ $=1$. The situation is a little messier when $d$ is not a prime power.)

Let $C[G]$ be the group-ring over the complex numbers of $G=$ $(z / d z) * /\{ \pm \pm\}$. Let $I$ be the ideal generated by the element $\sum_{\sigma \in G} \sigma$. Then it is easy to see that $E_{c y c}$ is a free rank-one $2[G] / 1$-module with generator $\varepsilon$, where we define

$$
\varepsilon^{\sum a_{\sigma} \sigma}=\prod(\sigma \varepsilon)^{a_{\sigma}}, \quad \sum a_{\sigma} \sigma \in Z[G] / L
$$

and $\mathrm{E}_{\mathrm{cyc}} \otimes \mathrm{C}$ is a free rank-one $\mathrm{C}[\mathrm{G}] / \mathrm{I}-$ module.
Let $X$ denote $C[G] / I$. Let $L O G ; X \longrightarrow X$ be the map defined on a basis element by

$$
\operatorname{LOG}\left(\sigma_{a}\right)=\sum_{b \in G} \log \mid \sigma_{b} \sigma_{a} \varepsilon!\sigma_{b}^{-1} .
$$

The determinant of LOG is clearly the regulator of $\mathrm{K}^{+}$(times the index $\left[E: E_{c y c}\right]$ ). We can write the map $L O G$ explicitly as

$$
\operatorname{LOG}\left(\sigma_{a}\right)=\sum_{b \in G} \log \left|\frac{\zeta^{a b g}-\zeta^{-a b g}}{\zeta^{a b}-\zeta^{-a b}}\right| \sigma_{b}^{-1} .
$$

It is easy to check that this is a well-defined map from $X$ to $X$.
An irreducible representation of $C[G] / I$ is the same as a nontrivial even character mod $d$. Let $X$ be such a character. Define the $\chi$-regulator $R_{\chi}$ to be the determinant of the map induced by LOG on $V_{X_{C[G]}}^{\otimes} X$, where $G$ acts on the 1-dimensional space $V_{X}$ by $X$. Here $V_{X_{C}} \otimes X$ can simply be identified with the $X$-eigen space of $X$, i.e., the 1 -dimensional subspace spanned by $\Sigma \bar{\chi}(a) \sigma_{a}$. Then

$$
\begin{aligned}
\operatorname{LOG}\left(\Sigma \bar{X}(a) \sigma_{a}\right) & =\sum_{a, b \in G} \bar{\chi}(a) \log \left|\sigma_{b} \sigma_{a} \varepsilon\right| \sigma_{b}^{-1} \\
& =\sum_{b, c \in G} X(b) \bar{\chi}(c) \log \left|\sigma_{c} \varepsilon\right| \sigma_{b}-1 \quad(c=a b) \\
& =\left(\sum_{a \in G} \bar{\chi}(a) \log \left|\sigma_{a} \varepsilon\right|\right)\left(\sum_{b \in G} \bar{\chi}(b) \sigma_{b}\right) .
\end{aligned}
$$

Thus

$$
\begin{aligned}
\mathrm{R}_{X} & =\sum_{a \in G} \bar{\chi}(a) \log \left|\frac{\zeta^{a g}-\zeta^{-a g}}{\zeta^{a}-\zeta^{-a}}\right| \\
& =\sum_{a \in G} \bar{\chi}(a)\left(\log \left|1-\zeta^{-2 a g}\right|-\log \left|1-\zeta^{-2 a}\right|\right) \\
& =\sum_{a \in G} \chi(2 g) \bar{\chi}(2 a g) \log \left|1-\zeta^{-2 a g}\right|-\sum_{a \in G} \chi(2) \bar{\chi}(2 a) \log \left|1-\zeta^{-2 a}\right| \\
& =\chi(2)(\chi(g)-1) \sum_{a \in G} \bar{\chi}(a) \log \left|1-\zeta^{-a}\right| \\
& =\frac{\chi(2)}{2}(\chi(g)-1) \sum_{a=1}^{d} \bar{\chi}(a) \log \left(1-\zeta^{-a}\right)
\end{aligned}
$$

(because $\chi$ is even, and $\left.\log \left(1-\zeta^{-a}\right)\left(1-\zeta^{a}\right)=\log \left|1-\zeta^{-a}\right|\left|1-\zeta^{a}\right|\right)$. Comparing with (1.1), we see that the only difference between $R_{\chi}$ and $L(1, \chi)$ is an algebraic factor; that is, $R_{X}$ is the transcendental part of $L(1, X)$. It is this fact which Stark and Tate generalize in their conjecture.

We get a companion fact about the behavior of $L(s, \chi)$ near $s=0$ if we use the functional equation for $L(s, X)$, which relates $L(s, X)$ to $L(1-s, \bar{X})$, and hence relates behavior near $s=1$ to
behavior near $s=0$. Suppose $X$ is a nontrivial even character. Then the functional equation is (see, e.g., [41], p. 5):

$$
L(s, X)=\frac{g_{X}}{2}\left(\frac{2 \pi}{d}\right)^{s} \frac{L(1-s, \bar{X})}{\Gamma(s) \cos (s / 2)} .
$$

If we let $s \longrightarrow 0$ and write $\Gamma(s)=\Gamma(s+1) / s$, we find that $L(0, \chi)=0$ (which we knew, since $B_{1, X}=0$ for $X$ even and nontrivial), and the Taylor expansion at $s=0$ starts out

$$
L(s, \chi)=s \frac{g_{X}}{2} L(1, \bar{X})+\text { higher terms. }
$$

Hence, the transcendental part of the first nonzero Taylor coefficient is the same as the transcendental part of $L(1, \bar{\chi})$, i.e., it is $R_{\bar{\chi}}$. Note that the non-vanishing of $R_{\bar{\chi}}$ implies non-vanishing of $L(1, \bar{X})$, and at 0 it implies that the zero of $L(s, X)$ is simple, i.e., $L^{\prime}(0, \chi) \neq 0$.

More generally, the Stark-Tate conjecture can be stated equivalently in terms of the behavior near either 1 or 0 , thanks to the functional equation.

In the p-adic case, there is no known (or expected) functional equation, and so there are two completely different p-adic analogs of the Stark-Tate conjecture, one at $s=1$ (due to Leopoldt and Serre), and one at $s=0$ (due to Gross).
2. Leopoldt's p-adic regulator

Let $K$ be a totally real number field, $n=[K: Q]$. By Dirichlet's unit theorem, the group $E$ of units of $K$ is the product of the roots of 1 in $K$ and a free abelian group of rank $n-1$. Let $e_{1}, \ldots, e_{n-1}$ be generators of this free abelian group. Let $\phi_{1}, \ldots$, $\phi_{\mathrm{n}}: \mathrm{K} \leftrightarrows \Omega_{\mathrm{p}}$ be all of the n possible imbeddings of K into the algebraically closed field $\Omega_{p}$. The Leopoldt (p-adic) regulator of $K$ is defined as the determinant of the $(n-1) \times(n-1)$ matrix
$\left\{\ln _{p} \phi_{i}\left(e_{j}\right)\right\}{ }_{1 \leq i, j \leq n-1}$.
Lemma. The Leopoldt regulator $R=R_{p, L e o p o 1 d t}{ }^{(K)}$ is independent up to $\pm 1$ of the choice of basis $\left\{_{j}\right.$, Leopoldt and the ordering of

## the $\phi_{i}$.

Proof. Any other basis $e^{\prime}=\left\{e_{j}^{r}\right\}$ can be written in the form $e^{\prime}=e^{M}$, where $M$ is an $(n-1) \times(n-1)$ matrix $\left\{m_{k j}\right\}$ with $m_{k j} \epsilon$ $z$ and $\operatorname{det} M= \pm 1$. The notation $e^{\prime}=e^{M}$ here means that $e_{j}^{\prime}={\underset{k}{i}}^{e_{k} m_{k j}}$. Then also $\phi_{i}\left(e_{j}^{\prime}\right)=\prod_{k} \phi_{i}\left(e_{k}^{m_{k j}}\right)$, and $\ln _{p} \phi_{i} e^{e^{\prime}}=$ ( $\left.\ln _{p} \phi_{i} e\right) M$, i.e., $\ln _{p} \phi_{i}\left(e_{j}^{\prime}\right)=\sum_{k} m_{k j} \ln _{p} \phi_{i}\left(e_{k}\right)$. This means that replacing $e$ by $e^{\prime}$ in (2.1) amounts to multiplying the matrix (2.1) by $M$ on the right. Since $\operatorname{det} M= \pm 1$, the independence of choice of basis $e$ is clear.

Rearranging the $\phi_{I}, \ldots, \phi_{\mathrm{n}-1}$ clearly changes the determinant at most by a sign. It remains to consider what happens if some $\phi_{\mathrm{k}}$, $1 \leq k \leq n-1$, changes places with $\phi_{n}$. Since $\prod_{i=1}^{n} \phi_{i}\left(e_{j}\right)=1$ for any unit $e_{j}$, we have $\sum_{i=1} \ln _{p} \phi_{i}\left(e_{j}\right)=0$, and so adding all the other rows to the $k$-th row gives $\quad \sum_{i=1} \ln _{p} \phi_{i}\left(e_{j}\right)=-\ln _{p_{n}} e_{j}\left(e_{j}\right)$ in the $k$-th row; hence interchanging $\phi_{k}$ and $\phi_{\mathrm{n}}$ only changes the determinant by a minus sign, and the lemma is proved.

The Leopoidt conjecture. $\mathrm{R}_{\mathrm{p}, \text { Leopoldt }}(\mathrm{K}) \neq 0$ for any totally real number field $K$.

In the simplest case, when $n=2$, i.e., $K$ is real quadratic, the non-vanishing of $R_{p}$, Leopoldt ${ }^{(K)}$ simply says that $1_{p}$ of a fundamental unit $e$ is nonzero. Since the kernel of the $\ln _{p}$ map consists of powers of $p$ times roots of unity, while ord $p$ of any unit is 0 and $e$ is not a root of 1 , it inmediately follows that $\ln _{\mathrm{p}} \mathrm{e} \neq 0$.

More generally, we shall prove Leopoldt's conjecture for all abelian extensions $K$ of $Q$. The proof relies upon the $p$-adic version of the following deep theorem of transcendence theory.

Baker's theorem [9]. If $0 \neq \alpha_{i} \in \bar{Q} \subset C$ and $\left\{\log \alpha_{i}\right\} \subset c$ are linearly independent over $Q$, then $\left\{\log \alpha_{i}\right\}$ are linea: $1 y$

## independent over $\bar{Q}$.

But before showing how the Leopoldt conjecture for $\mathrm{K} / \mathrm{Q}$ abelian can be derived from the p-adic Baker theoren, we first give another interpretation of this conjecture, in terms of which one can state a natural more general conjecture.

Let K be a number field with ring of integers 0 and group of units $E$. Let $[\mathrm{K}: \mathrm{Q}]=\mathrm{n}=\mathrm{r}_{1}+2 \mathrm{r}_{2}$, where $\mathrm{r}_{1}\left(2 \mathrm{r}_{2}\right)$ is the number of real (complex) imbeddings. Let $P_{i}, \mathbf{i}=1, \ldots, g$, be the primes of $K$ dividing $p$, and let $O_{p_{i}} \subset K_{p_{i}}$ be the $P_{i}$-adic completion of $0 \subset K$. Let $N_{i}$ denote the norm from $K_{P_{i}}$ to $Q_{p}$.

$$
A=\prod_{i=1}^{\text {Let }}{\stackrel{0}{P_{i}}}, \quad A_{0}=\left\{\left(x_{1}, \ldots, x_{g}\right) \in A \mid \Pi N_{i}\left(x_{i}\right)=1\right\}
$$

Let $E_{0} C E$ be the subgroup of units of norm $+1 ; E_{0}$ has index 2 in $E$ if 0 has units with norm -1 , otherwise $E_{0}=E$. Then $\mathrm{E}_{0} \subset 0 \hookrightarrow \mathrm{P}_{\mathrm{P}_{\mathrm{i}}}$ imbeds in $\mathrm{A}_{0}$.

Let $\bar{E}_{0}$ be the closure of $E_{0}$ in $A_{0}$. To get a concrete idea of what $\overline{\mathrm{E}}_{0}$ looks like, let $e_{1}, \ldots, e_{r_{1}+r_{2}-1} \in E$ be a set of fundamental units of norm +1 , i.e., they generate $E_{0}$ modulo roots of 1. Thus, $E_{0}=\left\{\eta \pi e_{j}^{\alpha_{j}} \mid \alpha_{j} \in Z, \eta\right.$ a root of 1$\}$. Now let $N$ be an integer such that $e_{j}^{N} \equiv 1\left(\bmod P_{i}\right)$ for all $i$ and j. For example, $N$ can be chosen to be $\prod_{i=1}^{g}\left(q_{i}-1\right)$, where $q_{i}=$ $p_{p}^{f_{i}}$ is the number of elements in the residue field $0 / P_{i}$. Let $e_{j}^{\prime}=e_{j}^{N}$. Then $e_{j}^{\prime}$ can be raised to p-adic powers $\alpha_{j} \in Z_{p}$ in $A_{0}$, because its image in each $0_{\hat{p}_{i}}$ is close to 1 . It is easy to see that $\bar{E}_{0} \subset A_{0}$ is precisely the set of elements of the form $n \pi e_{j}^{\beta_{j}} e_{j}^{e_{j}}$, where $\eta$ is a root of $1,0 \leq \beta_{j}<N$, and $\alpha_{j} \in Z_{p}$.

[^0]equivalent to the assertion that $\bar{E}_{0}$ is a subgroup of finite index in $A_{0}$.

Thus, if $K$ has $r_{2}$ pairs of complex imbeddings, a natural generalization of Leopoldt's conjecture is: $A_{0} / \bar{E}_{0}$ is isomorphic to (finite group) $\times z_{p}^{r_{2}}$.

Proof of proposition. First note that the fundamental units $e_{j}$ can be raised to some power $e_{j}^{\mathrm{Np}^{M}}$ such that the image of $e_{j}^{N_{p}^{M}}-1$ in $0_{P_{i}}$ has $p$-adic absolute value less than $p^{-1 /(p-1)}$ for all $i$, $j$. For example, if all of the $P_{i}$ are unramified (and $p>2$ ), then we need on1y take $e_{j}^{N}$, where $N$ is chosen as above so that $e_{j}^{N} \equiv 1\left(\bmod P_{i}\right) \quad$ for all $i, j$. Let $e_{j}^{\prime}=e_{j}^{N p^{M}}$, let $E^{\prime}=\left\{\Pi e_{j}^{\alpha_{j}} \mid\right.$ $\left.\alpha_{j} \in Z\right\} \subset E_{0}, \quad \bar{E}^{\prime}=\left\{\operatorname{Me}_{j}^{\prime} \alpha_{j} \mid \alpha_{j} \in Z_{p}\right\} \subset \bar{E}_{0}$. Then $E^{\prime}$ has finite index in $\mathrm{E}_{0}$, and $\overline{\mathrm{E}}^{\prime}$ has finite index in $\overline{\mathrm{E}}_{0}$.

Note that, if we replace $e_{j}$ by $e_{j}^{\top}$ in the definition of $R=$ $R_{p, \text { Leopoldt }}$, obtaining a new determinant $R^{\prime}$, the effect is to multiply the regulator by a nonzero constant (in fact, by ( $\left.\mathrm{Np}^{\mathrm{M}}\right)^{\mathrm{n}-1}$ or $2\left(\mathrm{~Np}^{M}\right)^{n-1}$, since each entry in the $(\mathrm{n}-1) \times(\mathrm{n}-1)$ matrix is multiplied by $\mathrm{N}^{\mathrm{M}}$, and we also have to throw in a 2 if $\mathrm{E}_{0} \neq \mathrm{E}$ ). Thus, the proposition is equivalent to: $R^{\prime}=0$ if and only if $\left[A_{0}: \bar{E}^{\prime}\right]<\infty$.

Let $n_{i}=\left[K_{i}: Q_{p}\right]$ be the local degree, and let $\sigma_{i t}, t=1$, $\ldots, n_{i}$, be the imbeddings $K_{P_{i}} \longrightarrow \Omega_{p}$. Let $K_{p_{i}, t}=\sigma_{i t}{ }_{K_{P_{i}}}$, $0_{P_{i}, t}=\sigma_{i t} 0_{P_{i}}$. Let $B=\Omega_{p}^{n}$, and let $B_{0}=\left\{\left(y_{1}, \ldots, y_{n}\right) \in B \mid\right.$ $\left.\sum y_{i}=0\right\}$, Define LOG: $A \rightarrow B$ by

$$
\begin{aligned}
& \operatorname{LOG}\left(x_{1}, \ldots, x_{g}\right)=\left(\ln _{p} \sigma_{11}\left(x_{1}\right), \ldots, \ln _{p} \sigma_{1, n_{1}}\left(x_{1}\right), \ldots,\right. \\
&\left.\ln _{p} \sigma_{g 1}\left(x_{g}\right), \ldots, \ln _{p} \sigma_{g \rightarrow n}\left(x_{g}\right)\right) .
\end{aligned}
$$

Let $O_{P_{i}}^{\prime}=\left\{x \in 0_{P_{i}}^{*}| | x-\left.1\right|_{p}<p^{-1 /(p-1)}\right\}$, which is a subgroup of finite index in $0_{\hat{P}_{i}}^{*}$. (Since $\left|\sigma_{i t}{ }^{x-1}\right|_{p}$ is independent of the imbadding, we denote it $|x-1|_{p}$.) Let $1+\pi_{i} \in 0_{P_{i}}^{\prime}$ be an element with $\left|\pi_{i}\right|_{p}$ maximal, and denote $\pi_{i t}=\sigma_{i t}\left(\pi_{i}\right)$. We claim that $1 n_{p}$ maps $O_{P_{i}, t}^{\prime}=\sigma_{i t} 0^{\prime} P_{i}^{\prime}$ isomorphically to $\pi_{i t} 0_{P_{i}}, t$. To see this, first check that for any $\alpha \in 0_{P_{i}}$,t the series

$$
\left(1+\pi_{i t}\right)^{\alpha}=\sum\binom{\alpha}{j} \pi_{i t}^{j}
$$

converges to an element of $O_{P_{1}, t}^{\prime}$ whose p-adic distance from 1 is $\left|\alpha \pi_{i t}\right|_{p}$. Since $l_{n_{p}}$ and the exponential function give mutually inverse isomorphisms between the open disc of radius $p^{-1 /(p-1)}$ around 1 and the open disc of radius $p^{-1 /(p-1)}$ around 0 , it follows that $0_{P_{i}}^{\prime}=\left(I+\pi_{i}\right){ }^{0_{P_{i}}} \xrightarrow{{ }^{1 n_{p} \circ \sigma_{t}}} \pi_{i t} 0_{P_{i}, t}$.

Thus, ${\underset{I}{i=1}}_{g}^{0_{P_{i}}^{\prime}} \quad$ is a subgroup of finite index in $g^{A}$ which is taken isomorphically by LOG to the free rank-one ${\underset{i=1}{1=1} p_{i}}_{p_{i}}$-module in $B$ generated by $\left(\ldots \operatorname{In}_{p}\left(\pi_{i t}\right) \ldots\right)_{i=1}, \ldots, g ; t=1, \ldots, n_{i}$. Since $\operatorname{rank}_{Z_{p}}+0_{P_{i}}=n$, it follows that $\operatorname{rank}_{Z_{p}} \operatorname{LOG}(A)=n$, and
 the $z_{p}$-submodule of $B_{0}$ spanned by

$$
\begin{equation*}
\left\{\left(\ldots \ln _{p} \sigma_{i t} e_{j}^{\prime} \cdots\right)_{i, t}\right\}_{j=1, \ldots, n-1}, \tag{2.2}
\end{equation*}
$$

it follows that

$$
\begin{aligned}
{\left[A_{0}: \bar{E}^{\prime}\right]<\infty } & \Leftrightarrow\left[\operatorname{LOG}\left(A_{0}\right): \operatorname{LOG}\left(\bar{E}^{\prime}\right)\right]<\infty \\
& \Leftrightarrow \operatorname{rank}_{Z_{p}} \operatorname{LOG}\left(\overline{\mathrm{E}}^{\prime}\right)=n-1 \\
& \Leftrightarrow \text { the set of vectors (2.2) has rank } n-1 \\
& \Leftrightarrow R^{\prime} \neq 0 . \quad \text { Q.E.D. }
\end{aligned}
$$

The proposition just proved can be paraphrased roughly as
follows: Leopoldt's conjecture says that a system of fundamental units is independent not only over $Z$, but even over $Z_{p}$.

We now prove Leopoldt's conjecture in the abelian over $Q$ case.
Theorem. Let $K$ be a totally real abelian extension of $Q$. Then $R_{p, \text { Leopoldt }}(K) \neq 0$.

Proof. Fix an imbedding $\phi: K \longrightarrow \Omega_{p}$. Let $G=G a 1(K / Q), \quad n=$ $[K: Q]$. Then the imbeddings $\phi_{i}: K \longrightarrow \Omega_{p}$ are $\left\{\phi^{\circ} \sigma\right\}_{\sigma \in G}$. Let $\sigma_{0}$ be any fixed element of $G$. Let $\left\{e_{j}\right\} \quad \begin{aligned} & \text { b } \\ & \text { be a basis of the units of }\end{aligned}$ K. Then

$$
R=R_{p, \text { Leopoldt }}(K)=\operatorname{Det}\left\{\ln _{p} \phi \sigma\left(e_{j}\right)\right\}{ }_{j=1}, \ldots, n-1 ; \sigma \in G-\left\{\sigma_{0}\right\}
$$

If $R=0$, then the rows of this matrix are linearly dependent over $\Omega_{p}$, i.e.,

$$
\sum_{\sigma \in G-\left\{\sigma_{0}\right\}} a_{\sigma}^{0} \ln _{p} \phi \sigma\left(e_{j}\right)=0, \quad j=1, \ldots, n-1,
$$

for some $a_{\sigma}^{0} \in \Omega_{p}$ not all zero. Since any unit $e$ is a root of 1 times a product of the $e_{j}$, we have
$\sum_{\sigma \in G} a_{\sigma}^{0} \ln p \sigma(e)=0, \quad a_{\sigma_{0}}^{0}=0, \quad$ for all units e .
Let $\Omega_{p}[G]$ be the group-ring over $\Omega_{p}$ of $G$. Define
$I=\left\{\sum_{\sigma \epsilon G}^{p} a_{\sigma} \sigma \in \Omega_{p}[G] \mid \sum_{\sigma \in G} a_{\sigma} \ln _{p} \phi \sigma(e)=0\right.$ for all units $\left.e\right\}$.
Then $I$ is an ideal of $\Omega_{p}[G]$, since it is clearly closed under addition, and for any $\tau \in G, \quad \sum a_{\sigma} \sigma \in I \Longrightarrow \sum_{\sigma} a_{\sigma} n_{p} \phi \sigma(\tau e)=0$ for all units $e \Rightarrow \tau \sum a_{\sigma} \sigma \in I$. By (2.3), $\sum a_{\sigma}^{0} \sigma$ is an element of $I$. Since $a_{\sigma_{0}}^{0}=0$, this element is not a multiple of $\sum \sigma$. Hence, we can find a nontrivial character $X: G \longrightarrow \Omega_{p}^{*}$ such that $\sum \mathrm{a}_{\sigma}^{0} \chi^{-1}(\sigma) \neq 0$. (This is because the function $\mathrm{f}(\sigma)=a_{\sigma}^{0}$ on $G$ can be expanded as a linear combination of characters of $G$ : $f=\sum_{X} c_{X} X$, with $c_{X}=\frac{1}{n} \sum_{\sigma} a_{\sigma}^{0} X^{-1}(\sigma) ;$ if $c_{X}=0$ for all nontrivial $X$, then $f$ would be a multiple of the trivial character,
$\Sigma \mathrm{a}_{\sigma}^{0} \sigma$ would be a multiple of $\Sigma \sigma$. )
So let $X$ be such that $\sum a_{\sigma}^{0} \chi^{-1}(\sigma) \neq 0$, and Iet
$\sigma_{X}=\sum X(\sigma) \sigma$.
Then, since $I$ is an ideal, it contains

$$
\sigma_{\chi} \sum a_{\sigma}^{0} \sigma=\sum_{\sigma} \sigma \sum_{\tau} a_{\tau}^{0} \chi\left(\tau^{-1} \sigma\right)=\left(\sum a_{\sigma}^{0} \chi^{-1}(\sigma)\right) \sigma_{\chi} .
$$

Since the coefficient is nonzero, it follows that $\sigma_{\chi} \in I$. Note that also $\sigma_{1}=\Sigma \sigma \in I$, because $\Pi \sigma(e)=1$ for all units e. Thus, $I$ contains $\sigma_{X}-\sigma_{1}=\sum_{\sigma \neq i d}(X(\sigma)-1) \sigma$, i.e.,
$\sum_{\sigma \neq i d}(\chi(\sigma)-1) \ln _{p} \phi \sigma(e)=0$ for all units e.
We now use the p-adic version of Baker's theorem, which was proved by Brumer [15]. It is the same theorem, except that log is replaced by ${1 n_{p}}$, $C$ is replaced by $\Omega_{p}$, and we fix an imbedding of $\bar{Q}$ in $\Omega_{p}^{p}$ instead of $c$. That is,
p-adic Baker theorem. If $0 \neq \alpha_{i} \in \bar{Q} \subset \Omega_{p}$ and $\left\{\ln _{p} \alpha_{i}\right\} \subset \Omega_{p}$ are linearly independent over $Q$, then they are linearly independent over $\bar{Q}$.

Because of this theorem, we may conclude from (2.5) that for all e the set $\left\{\ln _{p} \phi \sigma(e)\right\}_{\sigma \in G-\{i d\}}$ is linearly dependent over Q, i.e., over 2 . Thus, for every e there are integers $m_{\sigma}$ with $m_{i d}=0$ such that $\ln _{p} \phi\left(\Pi_{\sigma \in G} \sigma(e)^{m_{\sigma}}\right)=0$, i.e., $\phi\left(\operatorname{I\sigma }(e)^{\mathbb{m}_{\sigma}}\right)$ is a power of $p$ times a root of 1 . Since ord ${ }_{p}$ of any unit is zero, $\Pi \sigma(e){ }^{m_{\sigma}}$ must be a root of 1 . Thus, replacing $m_{\sigma}$ by a multiple, we obtain: for each e there exist $m_{\sigma}$ not all zero, but with $m_{i d}=0$, such that $\Pi \sigma(e)^{m_{\sigma}}=1$.

But, by a theorem of Minkowski ([73], p. 90), there exists a unit e such that $\Pi \sigma(e)^{\mathbb{m}_{\sigma}}=1, m_{i d}=0$, implies that all of the $\mathrm{m}_{\sigma}=0$. That is, there exist units whose conjugates are multipli-
catively independent except for the single relation $N e=\Pi \sigma(e)=1$. This contradiction proves the theorem.

The Leopoldt conjecture for all totally real fields would follow from the following conjecture in transcendence theory.

Conjecture (Schanuel). If $\alpha_{1}, \ldots, \alpha_{r} \in C$ are Inearly independent over $Q$, then

Tr. deg. ${ }_{Q} Q\left(\alpha_{1}, \ldots, \alpha_{r}, e^{\alpha_{1}}, \ldots, e^{\alpha_{r}}\right) \geq r$.
The same holds if $\alpha_{1}, \ldots, \alpha_{r} \in \Omega_{p}$ are in the disc of convergence of the $p$-adic exponential function and are linearly independent over Q.

To see how Leopoldt's conjecture would follow from Schanuel's conjecture, we shall suppose that $K$ is Galois (the general case can readily be reduced to the Galois case), in which case Minkowski's theorem cited above ensures the existence of a unit $e$ which together with its conjugates $\sigma_{i}(e), \sigma_{i} \in \operatorname{Gal}(K / Q)$, generates a subgroup of finite index in the unit group. Let $\phi_{i}=\phi \circ \sigma_{i}$ be the imbeddings $\mathrm{K} \leftrightarrows \Omega_{\mathrm{p}}$. In Schanuel's conjecture let $\mathrm{r}=\mathrm{n}-1, \alpha_{i}=$ $\operatorname{In}_{p} \phi_{i}(e), i=1, \ldots, n-1$. Replacing the full unit group by the subgroup generated by the $\phi_{i}(e)$ only changes the regulator by a nonzero constant multiple. If we set $e_{j}=\phi_{j}(e)$, we have $\phi_{i}\left(e_{j}\right)=$ $\phi\left(\sigma_{i} \sigma_{j}(e)\right)$. The regulator for the $e_{j}$ is then the determinant of a matrix each of whose rows is a permutation of

$$
\alpha_{1} \quad \alpha_{2} \cdots \alpha_{n-1} \quad\left(-\alpha_{1}-\alpha_{2}-\cdots-\alpha_{n-1}\right)
$$

with one entry omitted. Schanuel's conjecture says that $\alpha_{1}, \ldots$, $\alpha_{n-1}$ are algebraically independent. But vanishing of the regulator would give a nontrivial algebraic relation between the $\alpha$ 's. (The easiest way to see the non-triviality of the polynomial in $\alpha_{1}, \ldots, \alpha_{n-1}$ is to note that if it were the zero polynomial, then the classical regulator, which is the same determinant with $\alpha_{i}=$ $\log \left|\phi\left(\sigma_{i}(\mathrm{e})\right)\right|, \phi: \mathrm{K} \longrightarrow \mathrm{C}$, would also vanish, and it is well known that the classical regulator is nonzero.)
3. Gross's p-adic regulator

Let $K \subset C$ be a Galois extension of a totally real field $k$. Let $\tau$ be complex conjugation. Suppose we have an imbedding $\phi: \mathrm{K} \leftrightarrows \Omega_{\mathrm{p}}$ which extends a fixed imbedding $\mathrm{k} \leftrightarrows \Omega_{\mathrm{p}}$. Then any other such imbedding is of the form $\phi \circ \sigma, \sigma \in G a 1(\mathrm{~K} / \mathrm{k})$. By a complex conjugation of $\phi(K) \subset \Omega_{p}$ we mean the image of $\tau$ under any of these imbeddings $\phi \circ \sigma$, i.e., any of the automorphisms $(\phi \circ \sigma) \circ \tau \circ(\phi \circ \sigma)^{-1}=\phi \circ\left(\sigma \tau \sigma^{-1}\right) \circ \phi^{-1}$ of $\phi(\mathrm{K})$. If $\mathrm{k} / \mathrm{k}$ is abelian, then of course there is only one complex conjugation $\phi^{\circ} \tau^{\circ} \phi^{-1}$ of $\phi(\mathrm{K})$.

Let $k$ be totally real. Suppose we are given a representation $\rho: \operatorname{GaI}(\overline{\mathrm{k}} / \mathrm{k}) \longrightarrow \operatorname{Aut}(\mathrm{V})$,
where $V$ is a finite dimensional vector space over $\Omega_{p}$, which is trivial on $\operatorname{Gal}(\bar{k} / K) \subset G a 1(\bar{k} / k)$, i.e., $\rho$ can be considered as a representation of the quotient $\operatorname{Gal}(\mathrm{K} / \mathrm{k})$. If $\rho$ of any complex conjugation is the automorphism 1 (resp. -1), it is said to be an even (resp. odd) representation.

Using results of Deligne and Ribet [21], one can associate a $p$-adic L-function $L_{p}(s, \rho)$ to any even representation $\rho$. (If $\rho$ is not even, the associated $p$-adic L-function is identically zero.) $L_{p}(s, \rho)$ is a meromorphic function from $z_{p}$ to $\Omega_{p}$. It is conjectured to be holomorphic, except for a pole at $s_{s}=1$ when $\rho$ contains the trivial representation. $L_{p}(s, p)$ is called the "p-adic Artin L-series associated to $\rho$. "

Example. Let $k=Q$ and $K=Q(\zeta)$, where $\zeta$ is a primitive d-th root of $i$. Let $\operatorname{dim} V=1$, i.e., $\rho$ is a one-dimensional character. Such characters $\rho$ correspond to Dirichlet characters $X:(Z / d Z) * \longrightarrow \Omega_{\mathrm{P}}^{*}$ by the correspondence

$$
\rho\left(\sigma_{j}\right) e=\chi(j) e, \quad \sigma_{j} \in \operatorname{GaI}(K / Q) \approx(Z / d Z)^{*}
$$

where $e$ is a basis of $V=\Omega_{p} e$ and $j \in(Z / d Z)$ * is determined by $\sigma$ as usual by $\sigma_{j}(\zeta)=\zeta^{j}$. Then $\rho$ is even (resp, odd) if
$\chi(-1)=1$ (resp. $\chi(-1)=-1$ ). In this case the $p-a d i c$ Artin L-series associated to $\rho$ is simply the $p$-adic Dirichlet $L$-series $L_{p}(s, \chi)$ which we studied in Chapter II.

In this example, recall that $L_{p}(s, x)$ p-adically interpolates the algebraic numbers

$$
L_{p}(1-k, \chi)=\left(1-\chi \omega^{-k}(p) p^{k-1}\right) L\left(1-k, \chi \omega^{-k}\right),
$$

where $\omega$ is the Teichnüller character. Gross's conjecture concerns the expansion near $s=0$ of the $p$-adic Artin L-series $L_{p}(s, \rho)$. In the present example note that when $1-k=0$ in the above formula, the p-adic L-function is related to the classical L-function for the character $\chi \omega^{-1}$. Specifically,

$$
L_{p}(0, \chi)=\left(1-\chi \omega^{-1}(p)\right) L\left(0, \chi \omega^{-1}\right)=\left(1-\chi \omega^{-1}(p)\right)\left(-B_{1, \chi \omega^{-1}}\right)
$$

If $X$ is even, then $\chi \omega^{-1}$ is odd. The Gross conjecture for the expansion of $L_{p}(s, \rho)$ near $s=0$, which in some sense is a vast generalization of (3.1), will thus involve expressions associated to the odd representation $\rho \otimes \omega^{-1}$.

It will take us a while to work up to the precise statement of Gross's conjecture. We first define the "p-units" of a number field K:

$$
E=E^{(p)}(K) \underset{d e f}{=}\left\{\left.e \in K|\quad| e\right|_{v}=1 \text { for all valuations } v \|_{p}\right\}
$$

This means that (1) in the factorization of the fractional ideal (e) $=\operatorname{IP}^{\mathbb{m}_{p}}$ only $\left.P\right|_{P}$ occur; and (2) under all imbeddings $K \longrightarrow C$, e has complex absolute value $I$. Note that the p-units are not contained in the ring of integers $0 \subset \mathrm{~K}$. Condition (2) means that the $m$ for $P$ must be negative the $m$ for any complex conjugate prime ideal $\sigma \tau \sigma^{-1}(P)$. If all of the $m=0$, then it is well known that condition (2) implies that $e$ is a root of 1.

In the above example, when $K=Q(\zeta)$, an example of a p-unit is a ratio of Gauss sums of the form (see Chapter III)

$$
g\left(\tilde{X}_{a}, \psi_{\pi} \circ \operatorname{Tr}\right)^{d} / g\left(\tilde{X}_{a}^{-l}, \psi_{\pi} \circ \operatorname{Tr}\right)^{d}
$$

We shall see that these p-units play a crucial role in the case
when $K$ is abelian over $Q$, which is the one case where Gross's conjecture is proved.

The basic case which is of interest is when $K$ is a CM ("complex multiplication") field, i.e., a quadratic imaginary extension of a totally real field $\mathrm{K}^{+}$. In that case there is only one complex conjugation, namely, the unique nontrivial element $\tau$ of $\operatorname{Gal}\left(\mathrm{K} / \mathrm{K}^{+}\right)$, and we denote $\bar{\alpha}=\tau \alpha$ for $\alpha \in K$ and $\bar{P}=\tau P$ for a prime ideal $P$ of $K$. Thus, if $K$ is a $C M$ field, (1) and (2) give

$$
\begin{equation*}
e \in E^{(p)}(K) \quad \Longrightarrow \quad(e)=\left.\prod_{P}\right|_{p}(P / P)^{M_{P}} \tag{3.2}
\end{equation*}
$$

Writing (e) additively gives a homomorphism

$$
E^{(\mathrm{P})}(\mathrm{K}) \longrightarrow \oplus \mathrm{Z}(\mathrm{P}-\overline{\mathrm{P}})
$$

where the sum is over primes $P$ of $K$ dividing $p$, one from each complex conjugate pair. The kernel of this homomorphism is the group of roots of 1 in $K$, and the image certainly contains (-) $Z(P-\bar{P})$, where $h$ is the class number of $K$, because, if we write the principal ideal $P^{h}=(\alpha)$, we have

$$
\mathrm{E}^{(\mathrm{p})}(\mathrm{k}) \ni \alpha / \bar{\alpha} \longmapsto h(\mathrm{P}-\overline{\mathrm{P}}) .
$$

Thus, if we tensor the $Z$-module $E=E^{(P)}(\mathrm{K})$ (i.e., the abelian group with respect to multiplication, which we write additively) with $Q$ (thereby killing roots of 1 ), we obtain a Q-vector space

$$
\underset{Z}{E} \underset{Q}{\otimes} \approx \oplus Q(P-\tilde{P})
$$

We shall want to adjust the above homomorphism $\mathrm{E} \rightarrow(\mathrm{Z}(\mathrm{P}-\overline{\mathrm{P}})$ given by $e \longmapsto(e)=\Sigma m_{P}(P-\bar{P})$. Namely, at each $P$ insert the residue degree $f_{p}=\left[0 / P: F_{p}\right]$, where 0 is the ring of integers of $K$ and $F_{p}$ is the field of $p$ elements. Also, insert a minus sign. Thus, let
$\Psi(e) \underset{d e f}{=}-\sum m_{p} f_{p}(P-\bar{P})$.
This map $\Psi$ extends to an isomorphism $\Psi: E \otimes Q \xrightarrow{\sim} \oplus(P-\bar{P})$.
It is not hard to construct the inverse $\bar{\oplus}$ of the "divisor map" $\Psi$. Let $h$ be the class number of $K$, and write $p^{h}=(\alpha)$.

Let $e=\bar{\alpha} / \alpha \in E$. Then $e$ is determined by $P$. up to a root of unity, and so the element

$$
\begin{equation*}
\Phi(P-\bar{P})=\frac{1}{h f_{P}} e \in E \otimes Q \tag{3.3}
\end{equation*}
$$

is well-defined. Extend $\Phi$ by linearity to $\Theta \mathrm{Q}(\mathrm{P}-\overline{\mathrm{P}})$. Clearly $\Phi$ and $\psi$ are inverse to one another. These maps allow us to think of $E \otimes Q$ as divisors, and to think of any additive function on $E$ which kills roots of unity (for example, $\ln _{p}$ ) as a function on the divisors $\Theta \mathrm{Q}(\mathrm{P}-\overline{\mathrm{P}})$.

We now define a function LOG (not the same function as the LOG in the preceding section) by letting $\phi: K \hookrightarrow \Omega_{p}$ run through all imbeddings, letting $P_{\phi}$ be the prime ideal dividing $p$ which is defined by

$$
\begin{equation*}
P_{\phi}=\left\{\left.x \in 0|\quad| \phi(x)\right|_{p}<I\right\}, \tag{3.4}
\end{equation*}
$$

and setting

$$
\operatorname{LOG}(e)=\sum_{\phi} \operatorname{In}_{p} \phi(e) P_{\phi}, \quad e \in E^{(p)}(K) .
$$

Combining terms with the same $P_{\phi}$, we have

$$
\operatorname{LOG}(e)=\sum_{P} \operatorname{In}_{p}\left(N_{P}(e)\right)(P-\bar{P}) \in \Theta Q_{P}(P-\bar{P}),
$$

where the sum, as usual, is over $\left.P\right|_{P}$, one taken from each complex conjugate pair. Here $N_{p}(e)$ is the local norm $N_{K_{p} / Q_{p}}(\phi(e))$, where $\phi$ is any imbedding for which $P=P_{\phi}$. Since LOG kills roots of $I$ and is linear on $E$ (i.e., $\operatorname{LOG}\left(e_{1} e_{2}\right)=\operatorname{LOG}\left(e_{1}\right)+$ LOG( $e_{2}$ )), it extends uniquely to $E \otimes Q$, and so, via $\Phi$, to ( $+\mathrm{Q}(\mathrm{P}-\overline{\mathrm{P}})$ :

$$
\begin{equation*}
\text { LOG: } \oplus Q(P-\bar{P}) \longrightarrow \oplus Q_{p}(P-\bar{P}) \tag{3.5}
\end{equation*}
$$

Since LOG kills only roots of 1 in $E$, it is easy to see that its image in $\Theta Q_{p}(P-P)$ has $Q$-rank $g$, where $2 g$ is the number of primes $P$ over $p$.

But the interesting question is the $Q_{p}$-rank of the image. In other words, are the vectors $\operatorname{LOG}\left(e_{j}\right)$ even $Q_{p}$-independent as $\left\{e_{j}\right\}$ runs through a "fundamental set of p-units" (i.e., a maximal
set of p-units which are multiplicatively independent)? Gross conjectures that they are.

Gross's first conjecture. Let ${ }^{L} O Q_{Q}$ be the endomorphism of
$\oplus Q_{p}(P-\bar{P})$ obtained by extending $L O G^{p}$ linearly from $\oplus Q(P-\bar{P})$
to $\xlongequal{\mathrm{p}} \mathrm{Q}_{\mathrm{p}}(\mathrm{P}-\overline{\mathrm{P}})$. Define
$R_{p, \operatorname{Gross}}(K)=\operatorname{Det} \operatorname{LOG}_{Q}$.
Then

$$
R_{p, \operatorname{Gross}}(K) \neq 0
$$

We continue to let $K$ denote a $C M$ field, a purely imaginary quadratic extension of the totally real $f\left(y e 1 d K^{+}\right.$, and let $k$ denote the totally real ground field. Let $G=G a 1(K / k)$, and let $0: G \longrightarrow$ Aut(V) be a representation in a finite dimensional $\Omega_{\mathrm{p}}$ vector space $V$. (Note: Any continuous representation $\rho: \operatorname{Gal}(\overline{\mathrm{k}} / \mathrm{k})$ $\longrightarrow$ Aut(V) which is even (resp. odd) factors through Gal(K/k) for some CM field $K$, where $\rho$ is even (resp. odd) if its value on complex conjugation is 1 (resp. -1); since we shall on 1 y consider such $\rho$, there is no loss of generality in taking $K$ to be a C̣M field.)
$G$ also acts on the $\Omega_{\mathrm{p}}$-vector space
$x_{d e f}^{=} \oplus \Omega_{p}(P-\bar{P})$
by permuting the divisors $P-P$. We write $P^{\sigma}$ for $O P$; note that $\mathrm{P}^{\sigma \tau}=\sigma \mathrm{T} P=\left(\mathrm{P}^{\top}\right)^{\sigma}$. Note that complex conjugation acts by -1 We can combine the action of $G$ on $V$ with its action on $X$ by looking at the subspace

## $\left.(\mathrm{v} \otimes)_{\Omega_{p}}\right)^{G}$

of $G$-invariant elements in the tensor product. (In our discussion of $L(1, X)$ and the Stark-Tate conjecture in $\S 1$, we dealt with
$v_{\chi_{C}} \otimes_{G]} x$. For one-dimensional $x$, this is isomorphic to $\left(V_{\chi}-\otimes\right)^{G}$, reflecting the fact that the behavior of $L(s, \chi)$ near $s=1$ is
related by the functional equation to the behavior of $\mathrm{L}(\mathrm{s}, \bar{\chi})$ near $\mathrm{s}=0$. )

To see what $(V \otimes X)^{G}$ looks like, suppose that $G$ acts transitively on the primes $P$ of $K$ over $p$ (i.e., there is only one prime of $k$ over $p$, and let $P_{0}$ be a fixed prime ideal of $k$ over $p$. Let $D_{0} \subset G$ be the decomposition group of $P_{0}$. We shall suppose that the representation $\rho$ is odd. Then as a vector space $(V \otimes X)^{G}$ is isomorphic to the subspace $V^{D_{0}}$ of $V$ left fixed by $\rho\left(\mathrm{D}_{0}\right):$ namely, let $\mathrm{v} \in \mathrm{V}^{D_{0}}$ correspond to
$\sum_{\sigma \bmod D_{0}} \rho(\sigma) v \otimes P_{0}^{\sigma} \in(v \otimes x)^{G}$.
To give a simple example, suppose that $k=Q, K=Q(\zeta)$, $\zeta$ a primitive $d$-th root of $1, G \approx(Z / d Z) *$, and $\rho$ is one-dimensional. Thus, $\rho$ is given by a Dirichlet character $X: G \longrightarrow \Omega{ }_{P}$. Further suppose that $p=\Pi_{\sigma \in G} P_{0}^{\sigma}$ splits completely in $K$, i.e., $d \mid p-1$. In that case $(V \otimes)^{\sigma \in G}$ is spanned by the vector $\sum_{\sigma \in G} x(\sigma) P_{0}^{\sigma}$. (Compare with the definition (2.4) of $\sigma_{X} \in \Omega_{P}[G]$ in the proof of Leopoldt's conjecture for $K / Q$ abelian.) More generally, if $\mathrm{p} \not \equiv 1(\bmod d)$, then the same element spans $(v \otimes x)^{G}$ if $\chi(p)=1$, but $(V \otimes X)^{G}=0 \quad$ if $\quad x(p) \neq 1$.

Returning to the general case of an odd representation $\rho$ of $G=G a 1(K / k), \quad K$ a CM field and $k$ totally real, we see that the endomorphism $L O G_{\Omega_{p}}$ of $X=\bigoplus \Omega_{p}(P-\bar{P})$ is G-equivariant. This is because, for $e \in E=E^{(P)}(K)$ and $\sigma \in G$,
$\operatorname{LOG}(\sigma \mathrm{e})=\sum_{\phi} \ln _{\mathrm{p}} \phi(\sigma \mathrm{e}) \mathrm{P}_{\phi}=\sum_{\phi} \ln _{\mathrm{p}} \phi(\mathrm{e}) \mathrm{P}_{\phi \sigma}-1=\sum_{\phi} \ln _{\mathrm{p}} \phi(\mathrm{e}) \mathrm{P}_{\phi}^{\sigma}$,
since we have $P_{\phi \sigma^{-1}}=P_{\phi}^{\sigma}={ }_{\sigma} P_{\phi}$ directly from the definition (3.4).
Thus, $\operatorname{LOG}_{\Omega_{\mathrm{p}}}$ induces an endomorphism of $(V \otimes X)^{G}$, which we denote LOG $_{\mathrm{V}}$, Gross's regulator for $\rho$ is defined as

$$
R_{p, \operatorname{Gross}}(\rho)=\operatorname{Det} \operatorname{LoG}_{V}
$$

If the first Gross conjecture is true, then $\operatorname{LOG}_{V}$ is also an isomorphism, and $R_{p, G r o s s}(\rho) \neq 0$.

The earlier regulator $R_{p, G r o s s}(K)$ is a special case of $R_{p, \operatorname{Gross}}(\rho)$. Namely, let $k=K^{+}$, and let $\rho: \operatorname{GaI}\left(K / K^{+}\right) \longrightarrow\{ \pm 1\}$ be the unique nontrivial character. Then $(v \otimes x){ }^{G} \approx x$, and $R_{p, \operatorname{Gross}}(\rho)=R_{p, \text { Gross }}($ K $)$.

Recall that if $G=G a l(K / k)$ permutes the primes $P$ of $K$ lying over $p$--i.e., if there is only one prime of $k$ lying over $p$--then $\operatorname{dim}(v \otimes X)^{G}=\operatorname{dim} V^{D}$, where $D_{0}$ is one of the decomposition groups. More generally, if there is more than one prime of the ground field $k$ lying over $p$, we have the picture


Then for each $i$, $G$ permutes the $P_{i j}, j=1, \ldots, g_{i}$, lying over $P_{i}$, and the same argument shows that $(V \otimes X){ }^{G}$ is isomorphic as a vector space to $\oplus \mathrm{V}^{\mathrm{D}_{\mathrm{i}}}$, where $\mathrm{D}_{\mathrm{i}}$ is the decomposition group of $P_{i 1}$. We let
$m_{\rho}=\operatorname{dim}(V \otimes X)^{G}=\sum_{i} \operatorname{dim} V^{D} i$.
Then $R_{p, G r o s s}(\rho)$ is the determinant of an $m_{p} \times m_{\rho}$ matrix.
For example, if $V$ is one-dimensional, i.e., if $\rho: \operatorname{Gal}(\mathrm{K} / \mathrm{k})$ $\longrightarrow \Omega_{p}^{*}$ is a one-dimensional character, then $m_{\rho}$ is equal to the number of primes $P_{i}$ of $k$ lying over $p$ such that $\rho$ is trivial on one (and hence on all) of the decomposition groups $D_{i}$ of primes of $K$ lying over $P_{i}$. In the case of one-dimensional $\rho$, conjecturally the $m_{\rho}$ vanishing Euler factors $\left(1-\rho\left(P_{i}\right)\right)$ in the Deligne-Ribet [21] function $L_{p}(s, \rho \omega)$ at $s=0 \quad$ should lead to an $\mathbb{m}_{\rho}$-fold zero at $s=0$; and it is further conjectured [29] that the zero is of order exactly $m_{\rho}$. But it has not even been proved that $L_{p}(s, \rho \omega)$ has a multiple zero at $s=0$ when $m_{\rho}>1$.

Gross's conjecture, which presumes that the order of zero is at least $m_{\rho}$, concerns the coefficient of $s^{m_{\rho}}$ in the Taylor expansion at $s=0$ of $L_{p}(s, p \omega)$. In terms of his conjecture, we shall see that the assertion that $L_{p}(s, \rho \omega)$ has exactly an $m_{\rho}$-fold zero at $s=0$ is equivalent to non-vanishing of $R_{p, G r o s s}(\rho)$.

Gross's conjectured leading coefficient of $L_{p}(s, \rho \omega)$ at $s=0$ includes a certain algebraic number coming from the complex-analytic Artin L-series whose special values are p-adically interpolated by $L_{p}(s, p)$. We first recall the definition of the Artin L-series associated to a representation $\rho: \mathrm{Gal}(\mathrm{K} / \mathrm{k}) \longrightarrow$ Aut $(\mathrm{V})$, where now V is a finite dimensional complex vector space. Let $P$ be any prime of $k$ (not necessarily lying over $p$ ). Let $q=N_{k / Q} P$. Let $P$ be a prime of $K$ over $P$. Let $I_{p} \subset G a l(K / k)$ be the inertia group of $P$, and let $D_{P} \subset G a I(K / k)$ be its decomposition group. Let $F_{p} \in D_{p}$ be any automorphism such that $F_{p} \equiv \equiv x^{q}(\bmod P)$ for any $x$ in the ring of integers 0 of $K$. This Frobenius $F_{P}$ is uniquely deternined up to an element of $I_{p}$. Hence the "local factor at $P^{\prime \prime}$

$$
\operatorname{Det}\left(1-q^{-s} \rho\left(F_{p}\right) \mid V^{I_{p}}\right)
$$

where $V^{I_{p}}$ is the subspace of vectors fixed by $\rho\left(I_{p}\right)$, does not depend on the choice of Frobenius for $P$. If we change the choice of prime $P$ of $K$ lying over $P$, the effect is to conjugate $F_{P}$ and $I_{p}$ by some element of $\mathrm{Gal}(\mathrm{K} / \mathrm{k})$. Hence the determinant is unaffected. Thus, the above local factor depends only on $\rho, P$, and the complex variable $s$. The Artin L-series $L(s, p)$ is defined as the product of these local factors over all primes $P$ of $k$. This Euler product converges for $\operatorname{Re}(s)>1$ and has a meromorphic continuation onto the entire complex plane. The Artin conjecture asserts that it is holomorphic, except for a pole at $s=1$ if $\rho$ contains the trivial representation.

Example. Take the simple case when $k=Q, K=Q(\zeta), \quad \zeta$ a primitive d-th root of $1, G=G a l(K / k) \approx(Z / d Z) *$, and $\rho$ is a primitive character $X: G \rightarrow C^{*}$. If the ideal $(p)$ of $Q$ divides
d, then it ramifies in $K$, and $I_{p}$ (for any $P \mid p$ ) is the kerne 1 of the map $(z / d Z) * \longrightarrow\left(Z / d^{\prime} Z\right) *$, where $d^{\prime}=d / p^{\text {ord }}{ }_{p}{ }^{d}$. Since $\chi$ is primitive, it is nontrivial on $I_{p}$, and so $V^{I p}=0$. Thus, the local factor at $p$ is 1.

If, on the other hand, $p / d$, then $F_{p}=p \in(Z / d z) *$, and the local factor is $\left(1-p^{-s} X(p)\right)^{-1}$. Hence,

$$
L(s, p)=\prod_{p / d}\left(1-p^{-s} \chi(p)\right)^{-1}=\sum \frac{\chi(n)}{\mathfrak{n}^{s}}=L(s, \chi),
$$

which is the usual Dirichlet L-series.
Recall that in the case of Dirichlet L-series, before p-adically interpolating its values at negative integers we had to modify it by removing the Euler factor at $p$ :

$$
L^{*}(s, \chi)=\sum_{p / n} \frac{X(n)}{n^{s}}=\prod_{l \neq p}\left(1-\ell^{-s} X(\ell)\right)^{-1}=\left(1-p^{-s} X(p)\right) L(s, \chi) .
$$

A similar modification is required in the general case of Artin L-series before we can make the transition to p-adic Artin L-series. Namely, given our fixed prime $p$, we define the modified Artin L-series $L^{*}(s, \rho)$ to be the product of the local factors at all primes of $k$ not dividing $p$.
in the case of Dirichlet L-series, the values at negative integers $1-n$ are $-B_{n, X^{\prime}} / n \in Q(X)$, the field generated by the values of $\chi$. A similar fact was proved for Artin L-series $L(s, \rho)$ by Siegel [89]. Namely, first note that the representation $\rho: \mathrm{Gal}(\mathrm{K} / \mathrm{k})$ $\longrightarrow$ Aut(V) can be obtained by extension of scalars from a representation in a $K$-vector space $V_{K}$, where $K$ is a finite Galois extension of $Q$. (In other words, for a suitable basis of $V$, the matrix extries in $\rho(\sigma), \sigma \in \operatorname{Gal}(\mathrm{K} / \mathrm{k})$, are all in K.) Then Siegel showed that $L(1-n, \rho) \in K$, and, if $\sigma \rho, \sigma \in \operatorname{Gal}(K / Q)$, denotes the representation obtained by composing $\rho$ with the action of $\sigma$ on $\operatorname{Aut}\left(V_{K}\right)$, then $L(1-n, \sigma \rho)=\sigma L(1-n, \rho)$. The same is then clearly true of the modified L-function $\mathrm{L} *(1-\mathrm{n}, \mathrm{p})$.

By fixing once and for all an imbedding $\overline{\mathrm{Q}} \longrightarrow \Omega_{\mathrm{p}}$, we can
consider $\rho$ as giving a p-adic representation, which we also denote $\rho$, in $\mathrm{v}_{K} \otimes \Omega_{\mathrm{p}}$, which we also denote V . Then there exists a meromorphic function $L_{p}(s, \rho)$ on $Z_{p}$ with values in $\Omega_{p}$ which satisfies

$$
L_{p}(1-n, \rho)=L^{*}\left(1-n, \rho \otimes \omega^{-n}\right), \quad n \geq 2
$$

(where we use the fixed imbedding $\vec{Q} \longrightarrow \Omega_{p}$ to identify complex and $p$-adic representations and L-function values). When $n=1$, this relation

$$
\mathrm{L}_{\mathrm{p}}(0, \mathrm{p})=\mathrm{L} *\left(0, \rho \otimes \omega^{-1}\right)
$$

is also known to hold if $\rho$ is one-dimensional (or a direct sum of representations induced from one-dimensional representations); it is conjectured to hold for general $\rho$. $L_{p}(s, o)$ is identically zero unless $\rho$ is an even representation.

Gross now defines a second modification $L * *(s, \rho)$ of the Artin L-series. Recall that to get the first modification $L *(s, p)$, we threw out the local factors at primes $P$ of $k$ over $p$, by multiplying by the determinants (here $P$ is any prime of $K$ over P, $q=N P$ )
$\operatorname{Det}\left(1-q^{-s} \rho\left(F_{p}\right) \mid V^{I_{p}}\right)$.
To get $L^{\star *}(s, \rho)$, we put back in part of that local factor, by dividing by the subdeterminants

$$
\operatorname{Det}\left(1-q^{-s} \rho\left(F_{p}\right) \mid v^{D}\right)
$$

for $\left.P\right|_{P}$, where we restrict $F_{p}$ to the part of $V$ invariant under the whole decomposition group $\mathrm{D}_{\mathrm{P}}$. Since $\mathrm{F}_{\mathrm{P}} \in \mathrm{D}_{\mathrm{P}}$, we have $\rho\left(\mathrm{F}_{\mathrm{P}}\right) \mid \mathrm{V}^{\mathrm{D}} \mathrm{P}=$ identity, and so we define

$$
\begin{aligned}
& \mathrm{L} * *(\mathrm{~s}, \rho) \operatorname{dex} \overline{\mathrm{E}} \mathrm{~L}(\mathrm{~s}, \rho) \prod_{P \mid p}\left(1-q^{-s}\right)^{-\operatorname{dim} v^{D^{-}}} .
\end{aligned}
$$

For example, if $\rho$ is one-dimensional, this means that we put back in the Euler factors $\left(1-q^{-s} \rho\left(F_{p}\right)\right)^{-1}$ when $\rho\left(F_{P}\right)=1$. For instance, if $k=Q, K=Q(\zeta), \quad \zeta$ a primitive $d-t h$ root of 1 , $G=\operatorname{Gal}(K / k) \approx(Z / d Z) *$, and $\rho$ corresponds to the Dirichlet
character $X$, then $L^{*}(s, X)=\left(1-p^{-s} \chi(p)\right) L(s, \chi)$, and

$$
L * *(s, X)=\left\{\begin{array}{lll}
L *(s, \chi) & \text { if } & X(p) \neq 1 \\
L(s, \chi) & \text { if } & \chi(p)=1
\end{array}\right.
$$

The reason for this second modification is as follows. The subdeterminants $\operatorname{Det}\left(1-q^{-s} \rho\left(F_{p}\right) \mid V^{D_{P}}\right)$ of the factors Det $\left(1-q^{-s} \rho\left(F_{p}\right) \mid v^{T} p\right)$ that are thrown in to get $L^{*}(s, \rho)$ bring in zeros at $s=0$ of order $\operatorname{dim} V^{D} p$. Hence, $L^{*}(s, p)$ has an $m_{\rho}-$ fold zero at $s=0$, where $m_{\rho}=\sum_{P}$ dim $V{ }^{D} P$. Since $L_{p}\left(1-n, \rho \otimes \omega^{n}\right)$ interpolates the values $L_{i}^{*}(1-n, p)$, it is conjectured that $\mathrm{L}_{\mathrm{p}}(\mathrm{s}, \mathrm{p} \omega)$ also has an $m_{p}$-fold zero at $\mathrm{s}=0$, but this by no means follows from the mere fact that $L_{p}(s, \rho)$ interpolates $L^{*}(s, \rho)$.

To obtain the coefficient of the leading term of $L_{p}(s, \rho)$ at $s=0$, Gross therefore divides by the factors that conjecturally give the zeros at $s=0$.

Thus, the function whose value at $s=0$ is conjecturally related to this leading term is

$$
L^{* *}(s, \rho)=L(s, \rho) T \prod_{\operatorname{Det}}\left(1-q^{-s} \rho\left(F_{p}\right) \mid V^{I_{p}} / V^{D} p\right)
$$

where, as usual, the product is over all primes $P$ of $k$ over $p$, $q=N P$, and $P$ is some fixed choice of prime of $K$ over $P$ for each $P$.

Since $L_{p}(s, \rho)$ is only a nonzero function when $\rho$ is even, and since its value at $s=0$ is related to $L\left(s, \rho \otimes \omega^{-1}\right)$, if we replace $\rho$ by $\rho \otimes \omega$ we see that $L_{p}(s, \rho \otimes \omega)$ at $s=0$ should be related to $L(s, \rho)$, or rather $L * *(s, \rho)$, for $\rho$ an odd representation.

For an odd representation $\rho$, Gross defines (A stands for "algebraic part"):

$$
\mathrm{A}(\rho)=\mathrm{L} * *(0, \rho) \prod_{P} \mathrm{f}_{P}^{\operatorname{dim} \mathrm{V}^{\mathrm{D}} \mathrm{p}},
$$

where the product is over all primes $P$ of $k$ over $p, D_{p}$ is the
decomposition group of a prime $P$ of $K$ over $P$, and $f_{p}$ is the residue degree $\left[0 / P: F_{p}\right]$, where $o$ is the ring of integers of $k$ ( $f_{P}$ should not be confused with the residue degree [0/P: $F_{p}$ ] of $P$, or with the relative residue degree $[0 / P: 0 / P]=\# D_{P}$; in any case, this product term is just $l$ if $k=Q$.)

Without further ado, we can finally state Gross's main conjecture.

Gross's second conjecture. If $\rho: \mathrm{Gal}(\mathrm{K} / \mathrm{k}) \longrightarrow$ Aut(V) is an odd representation in a finite dimensional $\Omega_{\mathrm{p}}$-vector space v , then $L_{p}(s, p \otimes \omega)$ has a zero of order exactly $m_{\rho}=\Sigma d i m V^{D_{p}}$ at $s=0$, and

$$
\lim _{s \rightarrow 0} s^{-\mathbb{m}_{\rho}} L_{p}(s, \rho \otimes \omega)=R_{p, G \operatorname{coss}}(\rho) A(\rho)
$$

4. Gross's conjecture in the abelian over $Q$ case

We now prove this conjecture when $k=Q$ and $G=\operatorname{Gal}(\mathrm{K} / \mathrm{Q})$ is abelian. The conjecture is unproved in essentially any other case, even, for example, when $k=Q(\sqrt{D})$ and $K / k$ is abelian. Without loss of generality we may suppose that $K=Q(\zeta), \quad \zeta$ a primitive $d$-th root of 1 , since any abelian extension of $Q$ is contained in such a $K$, and all of the expressions in the conjecture remain the same if $K$ is replaced by a larger field. We first prove:

Proposition. Gross's first conjecture holds in the abelian over $Q$ case, i.e., $\operatorname{LOG}_{Q_{p}}$ is an automorphism of $\bigodot Q_{p}(P-P)$.

Proof. It suffices to prove that $\operatorname{LOG}_{\mathrm{Q}_{\mathrm{p}}}$ is an automorphism when we extend scalars from $Q_{p}$ to $\Omega_{p}$.
(We want to go to an algebraically closed field, so that we can decompose by the action of characters of $G=G a 1(K / Q) \approx(Z / d Z) *$.) Thus, we shall show that

$$
\begin{align*}
\text { LOG: } E \otimes \Omega_{p} & \longrightarrow \bigodot_{p}(P-\bar{P}) \\
e & \longmapsto \sum_{P} \ln _{p}\left(N_{p}(e)\right)(P-\bar{P}) \tag{4.1}
\end{align*}
$$

is an isomorphism, where the sum is over primes $P$ of $K$ over $P$,
one from each complex conjugate pair. (Recall that $N_{P}(e)$ denotes $\underset{\{x \in 0 \mid}{N_{K_{P}} / Q_{p}}(\phi(e))$, where $\phi$ is any imbedding $K \longrightarrow \Omega_{\mathrm{P}}$ for which $P=$ $\left\{\left.x \in 0|\quad| \phi(x)\right|_{p}<1\right\}$.)
$G$ acts on both $E \otimes \Omega_{p}$ and $\oplus \Omega_{p}(P-\bar{P})$, and we have seen that LOG is G-equivariant. Let us decompose both sides by characters $X$ of $G$. It is easy to see that the $X$-component of each side is at most one-dimensional; it is nonzero if and only if $\chi$ is odd and $x(p)=1$. In that case the $\chi$-component $\left(E \otimes \Omega_{p} \chi^{\chi}\right.$ is spanned by $e_{X \text { dé }} \Sigma \bar{\chi}(n) \sigma_{n}(e)$, where $e=\alpha / \bar{\alpha}$ is a $p$-unit with $(\alpha)=\cdot p^{h}$; and $\left(\oplus \Omega_{p}(P-\bar{P})\right)^{X}$ is spanned by $\Sigma \bar{X}(n) p^{\sigma_{n}}$, where $P$ is any fixed prime ideal of $K$ over $p$. It therefore suffices to show that $\operatorname{LOG}(\mathrm{e} \chi) \neq 0$.

If we denote $G=(Z / d Z) *, D=D_{P}=\left\{p^{j}\right\} \subset G$, and $f=\sharp D$, then by (4.1) we have

$$
\begin{aligned}
\operatorname{LOG}\left(e_{\chi}\right) & =2 £ \sum_{m \in G / \pm D} \sum_{n \in G / \pm D} \bar{\chi}(m) \ln _{p} N_{p}\left(\sigma_{n}\left(\sigma_{m}(e)\right)\left(P^{\sigma_{n}}-\bar{P}^{\sigma_{n}}\right)\right. \\
& =2 f \sum_{n, j \in G / \pm D} \bar{\chi}(n) \chi(j) \ln _{p} N_{p} \sigma_{j}(e)\left(P^{\sigma_{n}}-P^{\sigma_{n}}\right) \quad\left(j=\frac{n}{m}\right) \\
& =\left(\sum_{n \in G / \pm D} \chi(n) \ln _{p} N_{p} \sigma_{n}(e)\right)\left(\sum_{n \in G} \bar{\chi}(n)\left(p^{\sigma_{n}}-\bar{p}^{\sigma_{n}}\right)\right) .
\end{aligned}
$$

If this is zero, then from the p-adic Baker theorem (see §2) it follows that $1_{n}$ of the algebraic numbers $N_{p \sigma_{n}}$ (e) must be dependent over $Q$. For brevity let $e_{n}$ denote the product $\Pi \sigma_{n}^{-1}(\tau e)$ taken over all $\tau \in D$, so that $N_{p} \sigma_{n}(e)=N_{p}\left(\sigma_{n}^{-1}(e)\right)=$ $\phi\left(e_{n}\right)$ for $\phi$ any fixed imbedding $K \leftrightarrows \Omega_{p}$ for which $P=\{x \in 0\}$ $\left.|\phi(x)|_{p}<1\right\}$. Thus, for some $m_{n} \in Z$ not all zero we have

$$
\ln _{p}\left(\phi\left(\prod_{n \in G / \pm D} e_{n}^{m_{n}}\right)\right)=0
$$

This means that $\Pi e_{n}^{m_{n}}$ must be a power of $p$ times a root of unity; replacing $m_{n}$ by a suitable multiple, we obtain $\Pi{ }_{m} e_{n}^{m}=p^{r}$ for some $r \in Z$. But the ideal decomposition of $\Pi e_{n}^{m_{n}}{ }_{i n}^{n}$

$$
\left(\prod _ { n \in G / \pm D } \left(P_{n}^{\left.\left.\sigma_{n}^{-1} / \bar{p}^{\sigma_{n}}\right)^{m_{n}}\right)^{h f}, ~}\right.\right.
$$

while $(p)^{r}=\left(\prod_{n \in G / D} p^{\sigma_{n}}\right)^{r}$. This contradiction proves the proposition.

It is curious to note the resemblance between this proof and the proof of Leopoldt's conjecture for $K / Q$ abelian: in both cases the key step was to use the p-adic Baker theorem to conclude from the vanishing of a character sum that certain units are multiplicam tively dependent.

Theorem. Gross's second conjecture holds in the abelian over Q case.

Proof. Again $k=Q, K=Q(\zeta), \zeta$ a primitive d-th root of 1, $G=G a l(K / Q) \approx(Z / d Z) *$. Since any representation of the abelian group $G$ decomposes into a direct sum of one-dimensional characters and since both sides of Gross's conjecture can readily be verified to be multiplicative with respect to direct sums of representations, we can reduce to the case when $\rho$ is one-dimensional, i.e., is a Dirich1et character $\rho: G=(Z / \mathrm{dZ}) * \longrightarrow \overline{\mathrm{Q}}^{*}$. (We shall continue to use the letter $\rho$ for this Dirichlet character, rather than $X$, since the letter $X$ will soon be needed to denote a completely different kind of character, namely, a character of the multiplicative group of a finite field.)

Now there are two cases, depending on whether or not $\rho$ is trivial on the decomposition group $D=D_{P}=\left\{P^{j}\right\} \subset G \quad(P$ a prime of $K$ over $p$ ), i.e., we must consider: case (i) $\rho(p) \neq 1, m_{\rho}=$ $\operatorname{dim} V^{D}=0$; and case (ii) $\rho(p)=1, m_{\rho}=1$. Without loss of generality we may suppose that $d$ is the conductor of $\rho$; otherwise, $\rho$ factors through $\operatorname{Gal}\left(Q\left(\zeta^{\prime}\right) / Q\right)$, where $\zeta^{\top}$ is a primitive (cond $\rho$ ) -th root of 1 .

Case (i) $m_{p}=0$
In this case $R_{p, \operatorname{Gross}}(\rho)=1, A(\rho)=L^{*} *(0, \rho)=L^{*}(0, \rho)=$
$=(1-\rho(p)) L(0, \rho)$, and Gross's conjecture says that

$$
L_{p}(0, \rho \omega)=L(0, \rho)(1-\rho(p))
$$

which is true (see §II.4); in fact, both sides equal $-(1-\rho(p)) B_{1, \rho}$
We also know that this expression is nonzero (since $\rho$ is odd and primitive, and $\rho(p) \neq 1$ ), in other words, the order of vanishing
of $L_{p}(s, \rho w)$ at $s=0$ is $m_{\rho}=0$.
Case (ii) $m_{\rho}=1$, i.e., $\rho(p)=1$
This is the more interesting case.
We first compute $R_{p, G r o s s}(\rho)$, Let $\phi: K \longrightarrow \Omega_{p}$ be a fixed imbedding, so that $P=\left\{\left.x \in 0| | \phi(x)\right|_{p}<1\right\}$ denotes a fixed prime of K over $p$. As before, let $D=D_{p}=\left\{p^{j}\right\} \subset G=(z / d Z)$ *, and let $f=; D$. The one-dimensional vector space $(V \otimes X){ }^{G}$ is spanned by $\sum_{\sigma \in G / D} \rho(\sigma) P^{\sigma} \in X=\bigoplus_{\sigma \in G / \pm D} \Omega_{p}\left(P^{\sigma}-\bar{P}^{\sigma}\right)$.
Recall how LOG was defined on such an element. Let $\mathrm{P}^{\mathrm{h}}=(\alpha)$ Since $\left((\alpha / \bar{\alpha})^{\sigma}\right)=(P / \bar{P})^{\text {h }}$ for $\sigma \in \mathrm{D}$, it follows that $\alpha / \bar{\alpha}$ and $(\alpha / \bar{\alpha})^{\sigma}$ for $\sigma \in D$ differ by a root of unity (since their quotient is a unit of $K$ with complex absolute value 1). The image of the divisor $\sum_{\sigma \in G / D} \rho(\sigma) P^{\sigma}$ under the isomorphism $\Phi$ is (see (3.3))

$$
-\frac{1}{\mathrm{hf}} \sum_{\sigma \in G / \pm \mathrm{D}} \rho(\sigma)(\alpha / \alpha)^{\sigma} \in E \otimes \Omega_{\mathrm{p}},
$$

and, since

$$
\begin{aligned}
& \quad N_{\mathrm{p}^{\tau}}\left((\alpha / \bar{\alpha})^{\sigma}\right)=N_{\mathrm{R}_{\mathrm{P}}} / Q_{\mathrm{p}}\left(\phi \circ \tau^{-1}(\alpha / \bar{\alpha})^{\sigma}\right)=\phi\left(\left((\alpha / \bar{\alpha}) \tau^{\sigma \tau}\right)^{\mathrm{f}} \cdot \operatorname{root} \text { of } 1\right) \text {, } \\
& \text { it follows that }
\end{aligned}
$$

$$
\left.\begin{array}{rl}
\operatorname{LOG}\left(\sum_{\sigma \in G / D} \rho(\sigma) p^{\sigma}\right) & =-\frac{1}{h f} \sum_{\sigma \in G / \pm D} \rho(\sigma) \sum_{\tau \in G / D} 1 n_{p} N_{p}\left((\alpha / \tilde{\alpha})^{\sigma}\right) p^{\tau} \\
(\text { see }(3.5)
\end{array}\right)
$$

Hence

$$
R_{p, G r o s s}(\rho)=-\frac{1}{h} \sum_{\sigma \in G / D} \rho(\sigma) \ln _{p} \phi\left(\alpha^{\sigma}\right),
$$

where $\alpha$ is any generator of the ideal $P^{h}$.
Gross's conjecture asserts that

$$
L_{p}^{\prime}(0, \rho \omega)=R_{p, G r o s s}(\rho) A(\rho),
$$

where $A(\rho)=L(0, \rho)=-B_{1, \rho}$.
We now use the following results from earlier chapters:
(1) the formula for $L_{p}^{1}(0, \rho \omega)$ in $\S$ II. 8
$L_{p}^{\prime}\left(0, \rho(\omega)=\sum_{0<a<d} \rho(a) 1_{r_{p}} \Gamma_{p}^{\prime}(a / d)\right.$;
(2) the p-adic formula for Gauss sums in §III.6, which, after we take $\ln _{p}$ and use the fact that $\ln _{p} \pi=\frac{1}{p-1} \ln _{p}(-p)=0$, gives $\ln _{p} g\left(\tilde{X}_{a}^{-1}, \psi_{\pi}{ }^{\circ} T r\right)=\sum_{b \in D} \ln _{p} \Gamma_{p}(\langle a b / d\rangle)$
(see §III. 6 for notation; here $\tilde{X}_{a}$ is a multiplicative character of a finite field);
(3) Stickleberger's theorem (§III.7), which tells us that the ideal decomposition of $g\left(\tilde{x}_{a}^{-1}, \psi_{\pi}{ }^{\circ} \mathrm{Tr}\right)^{d}$, written additively, is $\sum_{j \in G}\langle a / j\rangle_{d} P^{\sigma_{j}}$, where $\left.<\right\rangle_{d}$ denotes least positive residue mod $d$.

$$
\begin{aligned}
& \text { We conclude that } \\
& L_{p}^{\prime}(0, \rho \omega)=\sum_{a \in G / D} \rho(a) \sum_{b \in D} \ln _{p} \Gamma_{p}(\langle a b / d\rangle) \\
& =\sum_{a \in G / D} \rho(a) \ln _{p} g\left(\tilde{X}_{a}^{-1}, \psi_{\pi}{ }^{\circ} T r\right) \\
& =\frac{1}{d h} \sum_{a \in G / D} \rho(a) \sum_{j \in G}\langle a / j\rangle{ }_{d} \operatorname{In}_{p} \phi\left(\alpha^{\sigma} j\right) \\
& \text { (where } \mathrm{P}^{\mathrm{h}}=(\alpha) \text { as before) } \\
& =\frac{1}{d h f} \sum_{a, j \in G} \rho(a / j)\langle a / j\rangle_{d} \rho(j) \ln _{p} \phi\left(\alpha^{\sigma} j\right) . \\
& =\frac{1}{d h f} \sum_{0<a<d} a \rho(a) \sum_{\sigma \in G} \rho(\sigma) \ln _{p} \phi\left(\alpha^{\sigma}\right) \\
& =\left(-\frac{1}{d} \sum_{0<a<d} a \rho(a)\right)\left(-\frac{1}{h} \sum_{\sigma \in G / D} \rho(\sigma) \ln _{p} \phi\left(\alpha^{\sigma}\right)\right) \text {, }
\end{aligned}
$$

which is precisely $\left(-B_{1, \rho}\right) R_{p, \operatorname{Gross}}(\rho)$, as desired. Q.E.D.
The above proof of Gross's first and second conjectures in the abelian over $Q$ case is Gross's variant of Ferrero-Greenberg's original proof in [29] of the simplicity of the zero of $L_{p}(s, \rho \omega)$ at $s=0$. The proof relies in an essential way upon several basic and diverse p-adic results: Brumer's p-adic Baker theorem, the formula relating $L_{p}^{\prime}(0, \rho \omega)$ to special values of the p-adic ganma function, Stickleberger's theorem, and the p-adic formula for Gauss sums. With this proof we conclude the main part of the book,

## APPENDIX

## 1. A theorem of Amice-Frèsnel

The measure $\mu_{z}$ on $z_{p}$ defined by $\mu_{z}\left(a+p^{N} Z_{p}\right)=z^{a} /\left(1-z^{p^{N}}\right)$, which was used to study p-adic Dirichlet L-functions in Chapter II, can also be used to give a simple proof of the following general fact.

Theorem (Amice and Frèsne1 [4]). Let $f(z)=\sum a_{n} z^{n} \in \Omega_{p}[\{z]]$ have the property that the coefficients $a_{n}$ can be $p$-adically interpolated, i.e., there exists a continuous function $\phi: \mathrm{z}_{\mathrm{p}} \Omega_{\mathrm{p}}$ such that $\phi(n)=a_{n}$. Then $£$ (whose disc of convergence must be the open unit disc $D_{0}\left(1^{\prime \prime}\right)$ is the restriction to $D_{0}\left(1^{-}\right)$of a Krasner analytic function $\tilde{f}$ on the complement of $D_{1}\left(I^{-}\right)$. In addition, $\tilde{f}$ has the Taylor series at infinity

$$
\tilde{f}(z)=-\sum_{n=1}^{\infty} \phi(-n) z^{-n}, \quad|z|_{p}>1
$$

Proof. Define
$\tilde{f}(z)=\int \phi d \mu_{z}$.
Then on $\Omega_{\mathrm{p}}-\mathrm{D}_{1}\left(1^{-}\right)$, the function $\tilde{f}$ is the uniform limit as $\mathrm{N} \longrightarrow \infty$ of the following rational functions with poles in $\mathrm{D}_{1}\left(1^{-}\right)$;

$$
\sum_{0 \leq n<p^{N}} \phi(n) \frac{z^{n}}{1-z^{p^{N}}}
$$

(see (2.3) of §II.2). Hence $\tilde{\mathbb{E}}$ is Krasner analytic on $\Omega_{p}-D_{1}\left(1^{-}\right)$, Next, for $|z|_{p}<I$ we have $\lim _{N \rightarrow \infty} 1 /\left(1-z^{p^{N}}\right) \doteq 1$, and so

$$
\tilde{\mathrm{f}}(\mathrm{z})=\lim _{\mathrm{N} \rightarrow \infty} \sum_{0 \leq \mathrm{n}<\mathrm{P}^{N}} \phi(n) z^{n}=f(z)
$$

Finally, if $|z|_{p}>1$, then we have

$$
\begin{aligned}
\sum_{0 \leq n<p^{N}} \phi(n) \frac{z^{n}}{1-z^{p^{N}}} & =\sum_{0<n \leq p^{N}} \phi\left(p^{N}-n\right) \frac{z^{-n}}{z^{-p^{N}}-1} \\
& \longrightarrow-\sum_{n=1}^{\infty} \phi(-n) z^{-n}
\end{aligned}
$$

as $N \longrightarrow \infty$, and the theorem is proved.
2. The classical Stieltjes transform

The Stieltjes transform of a function $f:[0, \infty) \longrightarrow C$ is
$G(z)=\int_{0}^{\infty} \frac{f(x)}{x+z} d x$
for all $z$ such that the integral converges. (More generally, one can replace $f(x) d x$ by $d F(x)$ and define the transform of that measure to be the corresponding Stieltjes integral.) Usually $f(x)$ is a rapidly decreasing function, and the integral converges for all $z \in C-(-\infty, 0]$. The Stieltjes transform is the square of the Laplace transform L:

$$
\begin{aligned}
L(L(f))(z) & =\int_{0}^{\infty} e^{-z t}\left(\int_{0}^{\infty} e^{-t x_{f}}(x) d x\right) d t \\
& =\int_{0}^{\infty} f(x) \int_{0}^{\infty} e^{-(z+x) t} d t d x=\int_{0}^{\infty} \frac{f(x)}{x+z} d x .
\end{aligned}
$$

It has been used extensively in the study of continued fraction expansions of analytic functions (this was Stieltjes' original purpose), numerical analysis, and quantum mechanics.

The function $£$ need not be rapidly decreasing in order for the Stieltjes transform (2.1) to exist. The Stieltjes transform
lso converges for $z \notin(-\infty, 0]$ if $f:[0, \infty) \longrightarrow C$ is a periodic function satisfying the conditions
(1) $f(x+d)=f(x)$;
(2) $f \in L^{1}([0, d])$;
(3) $\int_{0}^{d} f(x) d x=0$.

$$
\begin{aligned}
& \text { This is because } \\
& \qquad \begin{aligned}
|G(z)| & =\left|\sum_{n=1}^{\infty} \int_{(n-1) d}^{n d}\left(\frac{f(x)}{x+z}-\frac{f(x)}{n d}\right) d x\right| \text { by (2.4) } \\
& \leq \sum_{n=1}^{\infty} \frac{\text { const }}{(n d)^{2}} \int_{0}^{d}|f(x)| d x
\end{aligned}<\infty,
\end{aligned}
$$

where "const" depends on $z$ but not on $n$. It is periodic $f$ that we shall be particularly interested in from a number theoretic point of view.

We may suppose that in addition to (2.2)-(2.4) the function
$\mathrm{f}:[0, \infty) \longrightarrow \mathrm{C}$ satisfies
(4) $f(x)=0$ for $x<\delta$ for some positive $\delta$.

Example. Let $X$ be a nontrivial even (i.e., $X(-1)=1$ ) Dirichlet character of conductor $d$. Define

$$
\begin{equation*}
f_{x}(x)=\sum_{a=1}^{[x]} x(a) . \tag{2.6}
\end{equation*}
$$

Then $f_{X}$ obviously satisfies (2.2), (2.3) and (2.5) (with $\delta=1$ ).

$$
\begin{aligned}
& \text { To verify (2.4), we compute } \\
& \qquad \int_{0}^{d} f_{X}(x) d x=\sum_{a=1}^{d}(d-a) X(a)=-\sum_{a=1}^{d} a x(a)=-d_{1, X}=0
\end{aligned}
$$

for $X$ even. (If $X$ were odd, we would have to add the constant $B_{1, X}$ to $f_{X}$, and (2.5) would no longer hold; for simplicity in the discussion below, we want to assume (2.5).)

Suppose that
$\mathrm{f}:[0, \infty) \longrightarrow C$ satisfies (2.2)-(2.5). Then
(2.1) converges for $z \notin(-\infty,-\delta]$. For $|z|<\delta$ we expand $G(z)=\int_{0}^{\infty} \frac{f(x)}{x+z} d x$
$=\sum_{n=0}^{\infty}(-z)^{n} \int_{0}^{\infty} f(x) x^{-n-1} d x$.
Note that the "negative moments" $\int_{0}^{\infty} f(x) x^{-n-1} d x$ are convergent integrals, and are easily seen to be $\boldsymbol{O}\left(\delta^{-n}\right)$ as $n \longrightarrow \infty$. So (2.7) is the Taylor series of $G(z)$ in $|z|<\delta$.

$$
\begin{aligned}
& \text { In our example } f_{\chi}, \text { we compute for } s>0 \\
& \begin{aligned}
\int_{1}^{\infty} f_{\chi}(x) x^{-s-1} d x & =\sum_{k=1}^{\infty} \frac{1}{s}\left(\frac{1}{k^{s}}-\frac{1}{(k+1)^{s}}\right) \sum_{a=1}^{k} X(a) \\
& =\frac{1}{s} \sum_{k=1}^{\infty} \frac{\chi(k)}{k^{s}}=\frac{L(s, \chi)}{s}
\end{aligned}
\end{aligned}
$$

Thus for $f=f_{\chi}$ we have

$$
G(z)=\lim _{s \rightarrow 0} \frac{L(s, \alpha)}{s}+\sum_{n=1}^{\infty} \frac{L(n, \chi)}{n}(-z)^{n}, \quad|z|<1 .
$$

Returning to the general situation, suppose that $f:[0, \infty)$ $\longrightarrow C$ satisfies (2.2)-(2.4) (not necessarily (2.5)). Define $f^{(-1)}$ to be the integral of $f$ with constant of integration chosen so that $\int_{0}^{d} f^{(-1)}(x) d x=0$; thus, if $f$ is represented by the Fourier series $\sum_{a_{n}} e^{2 \pi i n x / d}$, with $a_{0}=0$ because of (2.4), then $f^{(-1)}$ has Fourier series $\frac{d}{2 \pi i} \sum_{n \neq 0} \frac{a_{n}}{n} e^{2 \pi i n x / d}$. Then define $f^{(-j-1)}$ inductively as $\left(f^{(-j)}\right)^{(-1)}, j=1,2, \ldots$.

We obtain an asymptotic series for $G(z)$ as $z \longrightarrow \infty$ (along any ray other than the negative real axis) by integrating by parts:

$$
\begin{aligned}
& G(z)= \int_{0}^{\infty} \frac{f(x)}{x+z} d x= \\
&=-\frac{f^{(-1)}(0)}{z}+\int_{0}^{\infty} \frac{f^{(-1)}(x)}{(x+z)^{2}} d x=\cdots \\
& z-\frac{f^{(-2)}(0)}{z^{2}}-\cdots-\frac{(j-1)!f^{(-j)}(0)}{z^{j}}-\cdots \\
&-\frac{(n-1)!f^{(-n)}(0)}{z^{n}}+n!\int_{0}^{\infty} \frac{f^{(-n)}(x)}{(x+z)^{n+1}} d x
\end{aligned}
$$

for any $n$. For fixed $n$, if $z$ approaches infinity away from the negative real axis, then it is easy to see that the integral term is $\boldsymbol{O}\left(\mathrm{z}^{-\mathrm{n}-1}\right)$. Thus, $-\sum(\mathrm{n}-1)!\mathrm{f}^{(-\mathrm{n})}(0) \mathrm{z}^{-\mathrm{n}}$ is an asymptotic series for $G(z)$.

For example, if $f=f_{\chi} \quad$ ( $X$ a nontrivial even character), we

$$
\begin{aligned}
& \text { have the easily computed Fourier series expansion } \\
& f_{X}(x)=\sum_{a=1}^{[x]} x(a) \\
& =\sum_{n \neq 0} a_{n} e^{2 \pi i n z / d} \text {, where } a_{n}=\frac{g_{\chi}}{2 \pi i} \frac{\bar{\chi}(n)}{n} \\
& \text { (here } g_{X}=\sum_{0<a<d} X(a) e^{2 \pi i a / d} \text { is the Gauss sum for } \chi \text { ). Then } \\
& f^{(-j)}(0)=\left(\frac{d}{2 \pi i}\right)^{j} \sum_{n \neq 0} \frac{a_{n}}{n^{j}}=\frac{d^{j} g_{X}}{(2 \pi i)^{j+1}} \sum_{n \neq 0} \frac{\bar{x}(n)}{n^{j+1}} \\
& = \begin{cases}\frac{d^{j} g_{X}}{(2 \pi i))^{j+I}} 2 L(j+1, \bar{\chi}), & j \text { odd } ; \\
0, & j \text { even }\end{cases} \\
& =\frac{1}{j!} L(-j, \chi)
\end{aligned}
$$

by the functional equation for $L(s, X)$ (see, e.g., [41], p. 104).

Thus, for $f=f_{\chi}$ we have

$$
\begin{equation*}
G(z) \sim-\sum_{j=1}^{n} \frac{L(-j, x)}{j} z^{-j}=-\sum_{j=1}^{n} \frac{L(-j, x)}{-j}(-z)^{-j} \tag{2.9}
\end{equation*}
$$

(since $L(-j, \chi)=0$ for $j$ even).
Comparing (2.8) and (2.9), we see that we have established a special case of the following theorem of Mellin-LeRoy (see [65], p. 109, 113).

Theorem. Suppose $\phi(s)$ is a function which is hoiomorphic and bounded on $\operatorname{Re} s \geq 0$. Then the Taylor series $G(z)=\sum \phi(n) z^{n}$ extends holomorphically onto $C-[0, \infty)$ and has asymptotic series $G(z) \sim-\Sigma \phi(\sim n) z^{-n}$ as $z \longrightarrow \infty$ (along any ray other than the positive real axis).

In our example $\phi(s)=\frac{1}{s} L(s, \chi) \quad$ (which is bounded even as $s \longrightarrow 0$, since $L(0, \chi)=0$ for $X$ even and nontrivial), and we have replaced $z$ by $-z$.

Remarks. 1. This theorem is the classical analog of the p-adic theorem of Amice-Frèsnel in $\S 1$.
2. Under a weaker assumption on $\phi(s)$, namely, bounded exponential growth, one has the same conclusion, except that not only $[0, \infty)$ but a whole sector $|\operatorname{Arg} z| \leq \theta$ must be excluded from the region where $G(z)$ is defined.
3. When $f=f_{\chi}, G(z)$ is a "twisted" (by $X$ ) $\log$ gamma
ction. In fact, we have function. In fact, we have

Proposition. Let $x$ be a nontrivial even Dirich1et character, and let $f$ be defined by (2.6). Then the Stieltjes transform $G$ of $\int_{0}^{\infty} \frac{i s}{f^{\frac{i s}{(x)}}} \frac{x+z}{x+z}=\sum_{a=1}^{d} x(a) \log \Gamma\left(\frac{z+a}{d}\right)$.

Proof. Let $G(z)$ be the left side of (2.10). Then

$$
\begin{aligned}
G(z) & =\lim _{n \rightarrow \infty} \int_{0}^{n d} \frac{f(x)}{x+z} d x \\
& =\lim _{n \rightarrow \infty} \sum_{j=0}^{d n-1}\left((\log (z+j+1)-\log (z+j)) \sum_{a=1}^{j} x(a)\right) \\
& =\lim _{n \rightarrow \infty}-\sum_{j=0}^{d n-1} \log (z+j) x(j) \\
& =\lim _{n \rightarrow \infty}-\sum_{a=1}^{d} x(a) \sum_{j=0}^{n-1} \log (z+d j+a) \\
& =\lim _{n \rightarrow \infty}-\sum_{a=1}^{d} x(a) \sum_{j=0}^{n-1} \log \left(\frac{z+a}{d}+j\right) .
\end{aligned}
$$

On the other hand, using the standard formula

$$
\Gamma(z)=\lim _{n \rightarrow \infty} \frac{(n-1)!n^{z}}{z(z+1) \cdots(z+n-1)},
$$

we see that the right side of (2.10) equals

$$
\lim _{n \rightarrow \infty} \sum_{a=1}^{d} x(a)\left(\log (n-1)!+\frac{z+a}{d} \log n-\sum_{j=0}^{n-1} \log \left(\frac{z+a}{d}+j\right)\right),
$$

and the first two terms vanish, because $\Sigma \chi(a)=\Sigma a \chi(a)=0$ for $X$ nontrivial and even. This proves (2.10).
4. If $X$ is odd or the trivial character, then, in addition to the Stieltjes transform $G(z)$, the asymptotic series for the twisted 10 g gamma function on the right in (2.10) also includes a principal term. For example, in the case of the trivial character, we have the "Stirling series" (see [97], p. 261)

$$
\log \frac{\Gamma(z)}{\sqrt{2 \pi}}=\left(z-\frac{1}{2}\right) \log z-z-\int_{0}^{\infty} \frac{x-[x]-1 / 2}{x+z} d x
$$

Note that for the trivial character we take $f_{X_{\text {triv }}}(x)=\sum \frac{e^{2 \pi i n x}}{2 \pi i n}$ $=-B_{1}(x)=[x]-x+\frac{1}{2}$. In other words, the Stieltjes transform of
the first Bernoulli polynomial gives the error in Stirling's formula $n!\sim \sqrt{2 \pi n} \frac{n^{n}}{e^{n}}$.
5. As mentioned before, the classical Stieltjes transform can be defined more generally for a Stieltjes measure $\mu=\mathrm{dF}(\mathrm{x})$ on the positive reals. For example, if $X$ is a nontrivial even Dirichlet character of conductor $d$, then the derivative of

$$
\sum_{a=1}^{d} x(a) \log \Gamma\left(\frac{z+a}{d^{v}}\right)=-\lim _{n \rightarrow \infty} \sum_{j=0}^{d n-1} x(j) \log (z+j)
$$

(see the proof of (2.10)) is

$$
-\lim _{n \rightarrow \infty} \sum_{j=0}^{d n-1} \frac{x(j)}{z+j}=-\int_{0}^{\infty} \frac{d f^{\prime}(x)}{x+z}
$$

where ${ }^{f} \chi$ is the function (2.6), i.e., $d f_{X}$ has point mass $\chi(j)$ at $j$. Thus,

$$
\begin{equation*}
\frac{d}{d z} \sum_{a=1}^{d} \chi(a) \log \Gamma\left(\frac{z+a}{d}\right)=-\int_{0}^{\infty} \frac{d f_{\chi}(x)}{x+z} . \tag{2.11}
\end{equation*}
$$

This formula can also be obtained by differentiating (2.10) under the integral sign and then integrating by parts:
$\frac{d}{d z} \sum_{a=1}^{d} X(a) \log \Gamma\left(\frac{z+a}{d}\right)=-\int_{0}^{\infty} \frac{f_{X}(x)}{(x+z)^{2}} d x$

$$
=\left.\frac{f_{\chi}(x)}{x+z}\right|_{0} ^{\infty}-\int_{0}^{\infty} \frac{d f_{\chi}(x)}{x+z}=-\int_{0}^{\infty} \frac{d f_{\chi}(x)}{x+z} .
$$

The formula (2.11) is closely analogous to the formula for the derivative of the twisted p-adic log gamma function (see (8.6) in Chapter II). More precisely, define the p-adic log gama function twisted by a nontrivial even character $X$ as follows:

$$
G_{p, X}(z) d={ }_{d} \sum_{a=I}^{d} X(a) G_{p}\left(\frac{z+a}{d}\right)
$$

$=\lim _{n \rightarrow \infty} p^{-n} \sum_{0<a<d, 0<j<p^{n}} x(a)\left(\frac{z+a}{d}+j\right)\left(\ln _{p}\left(\frac{z+a}{d}+j\right)-1\right)$
$=\lim _{\mathrm{n} \rightarrow \infty} \frac{1}{{d p^{n}}^{n}} \sum_{0<j<d p^{n}} x(j)(z+j)\left(\operatorname{In}_{p}(z+j)-i\right)$.
Then $G_{p, X}$ can be expressed in terms of $G_{p, \xi^{(z)}}=\lim _{n \rightarrow \infty} \frac{1}{d p^{n}}$ $\sum_{0 \leq j^{<} p^{n}} \xi^{j}(z+j)\left(\ln _{p}(z+j)-1\right)$, where $\xi^{d}=1, \quad \xi \neq 1, \quad$ as follows: $G_{p, \chi^{(z)}}=\frac{g_{X}}{d} \sum_{b=1}^{d} \bar{X}(b) G_{p, \xi^{b}}^{(z),}$
where $\xi$ is a fixed primitive d-th root of unity and $g_{X}=\sum_{a=1}^{d}$ $\chi(a) \xi^{a}=d / g_{\chi}$. So if we define a measure $\mu_{\chi}$ on $Z_{p}$ by

$$
\begin{equation*}
\mu_{X}=\frac{g_{X}}{d} \sum_{b=1}^{d} \bar{X}^{\mathrm{X}}(\mathrm{~b}) \mu_{\xi^{b}} \text {, where } \mu_{\xi^{b}}\left(a+p^{N} Z_{p}\right)=\frac{\xi^{a b}}{1-\xi^{\mathrm{b} p^{\mathrm{N}}}} \tag{2.12}
\end{equation*}
$$

$$
\begin{aligned}
& \text { then we have } \\
& \qquad G_{p}, X^{(z)}=-\int_{Z_{p}} \ln _{p}(x+z) d \mu_{X}(x) ; \\
& \frac{d}{d z} \sum_{a=1}^{d} X(a) G_{p}\left(\frac{z+a}{d}\right)=-\int_{Z_{p}} \frac{d \mu_{X}(x)}{x+z},
\end{aligned}
$$

which is the p-adic analog of ${ }^{p}(2.11)$.
Final remark. In the classical case an integer a prime to
d has point mass
$\quad \mathrm{df}_{\chi}(\mathrm{a})=\mathrm{g}(\mathrm{a})=\frac{g_{X}}{\mathrm{~d}} \sum_{\mathrm{b}=1}^{\mathrm{d}} \bar{\chi}(\mathrm{b}) \xi^{\mathrm{ab}}$.
Compare (2.13) to (2.12). As in our discussion of Leopoldt's formula for $L_{p}(1, \chi)$ in $\S I I .5$, we see that the $p$-adic construction is formally analogous to the classical case inside the open unit
 but the p-adic case only becomes arithmetically interesting when we extend to roots of unity, which are on the "boundary" of the unit disc.

In the remaining sections we shall give a systematic account of
the p-adic Stieltjes transform, following Vishik [95]. First we introduce a type of p-adic integration (not to be confused with the type used in Chapter II) which is the tool used to construct the inverse Stieltjes transform in the p-adic case.
3. The Shnirelman integral and the p-adic stieltjes transform

A p-adic analog of the line integral was introduced by Shnirelman in 1938 [88]. It can be used to prove p-adic analogs of the Cauchy integral theorem, the residue theorem, and the maximum modulus principle of complex analysis. The main applications of the Shnirelman integral are in transcendental number theory (see [1], [17]). Our interest in it will be to construct the inverse Stieltjes transform.

Definition. Let $f(x)$ be an $\Omega_{p}$-valued function defined on al1 $x \in \Omega_{p}$ such that $|x-a|_{p}=r$, where $a \in \Omega_{p}$ and $r$ is a positive real number. (We shall always assume that $r$ is in $\left|\Omega_{p}\right|_{p}$, i.e., a rational power of $\left.p.\right)$ Let $T \in \Omega_{p}$ be such that $|\Gamma|_{p}=r$. Then the Shnirelman integral is defined as the following limit if it exists:

$$
\int_{a, \Gamma} f(x) d x \operatorname{def}^{=} \underset{n \rightarrow \infty}{1 i m^{\Gamma}} \frac{1}{n} \sum_{\xi^{\mathrm{n}}=1} f(a+\xi \Gamma)
$$

where the ' indicates that the limit is only over $n$ not divisible by $p$.

Lemma 1. (1) If $\int_{a, \Gamma} f(x) d x$ exists, then
$\left|\int_{a, \Gamma} f(x) d x\right|_{p} \leq \max _{|x-a|_{p}=r}|f(x)|_{p}$.
(2) $\int_{a, \Gamma}$ commutes with limits of functions which are uniform limits on $\left\{x\left||x-a|_{p}=r\right\}\right.$.
(3) If $r_{1} \leq r \leq r_{2}$ and $f(x)$ is given by a convergent Laurent
$\frac{\text { series }}{\int f(x)} \sum_{k=-\infty}^{\infty} c_{k}(x-a)^{k}$ in the annuius $r_{1} \leq|x-a|_{p} \leq r_{2}$, then
$\int_{a, \Gamma} f(x) d x$ exists and equals $c_{0}$. In particular, the integral does a, $\Gamma$ depend on the choice of $\Gamma$ with $|\Gamma|_{p}=r$ or even on $r$, as long, as $r_{1} \leq r \leq r_{2}$. More generally,

$$
\int_{a, \Gamma} f(x)(x-a)^{-k} d x=c_{k} .
$$

The proof of the lemma is easy. Part (3) uses the fact that if $k \neq 0$, then $\sum_{\xi^{n}=1} \xi^{k}=0$ for $n>|k|$.

Lemma 2. For fixed $z \in \Omega_{p}$ and for $m>0$

$$
\int_{a, \Gamma} \frac{d x}{(x-z)^{m}}= \begin{cases}0 & \text { if }|z-a|_{p}<r \\ (a-z)^{-m} & \text { if }|z-a|_{p}>r\end{cases}
$$

To prove this, note that for $|x-a|_{p}=r$ we have the Laurent expansion

$$
\frac{1}{(x-z)^{m}}=\left\{\begin{array}{lll}
\left(\sum_{k=0}^{\infty}(z-a)^{k}(x-a)^{-k-1}\right)^{m} & \text { if } & |z-a|_{p}<r  \tag{3.1}\\
\left(\frac{1}{a-z} \sum_{k=0}^{\infty}(z-a)^{-k}(x-a)^{k}\right)^{m} & \text { if } & |z-a|_{p}>r
\end{array}\right.
$$

Then use part (3) of Lemma 1 (with $r_{1}=r_{2}=r$ ).
Lemma 3. (1) If $f(x)$ is a function on the closed disc of radius $\underset{p}{r}$ with center a, i.e., $f: D_{a}(r) \longrightarrow \Omega_{p}$, and if $f(x)=\sum_{k=0}^{\infty} c_{k}(x-a)^{k}$ with $r^{k}\left|c_{k}\right|_{p} \longrightarrow 0$, define $\|f\|_{r}=\max _{k}$ $r^{k}\left|c_{k}\right|_{p}$. Then $\max _{x \in D_{a}(r)}|f(x)|_{p}$ is attained when $|x-a|_{p}=r$ and equals $\|f\|_{r}$.
(2) Any Krasner analytic function $f: D_{a}(r) \longrightarrow \Omega_{p}$ (i.e., $f$ is a uniform 1imit of rational functions with poles outside $D_{a}(r)$ ) is of the form in part (1), i.e., is given by a power series.

Proof. Making a linear change of variables, without loss of generality we may assume that $a=0$ and $r=1$. Multiplying $f$ by a constant, we may also suppose that $\|f\|_{1}=\max \left|c_{k}\right|_{p}=1$. Clearly $|f(x)|_{p}=\left|\Sigma c_{k} x^{k}\right| \leq 1$ for $x \in D_{0}(1)$. Let $\bar{f}(x)=\Sigma \bar{c}_{k} x^{k} \in \bar{F}_{p}[x]$ be the reduction modulo $M_{\Omega_{p}}$, the maximal ideal in $\Omega_{p}$. If $x \in$ $D_{0}(1)$ is any element whose reduction mod $M_{\Omega_{p}}$ is nonzero and is not a root of the polynomial $\bar{f}$, then $|x|_{p}^{p}=1$ and $|f(x)|_{p}=1$. This proves part (1)

It follows from part (1) that, if $f_{n}$ is a sequence of rational functions approaching $f$ uniformly, and if each $f_{n}$ is represented by a power series on $D_{a}(r)$, then the sequence of power series approaches (coefficient by coefficient) a power series which represents $f$ on $D_{a}(r)$. Thus, it suffices to consider the case when $f$ is a rational function with poles outside $D_{a}(r)$. Again we make a change of variables so that $a=0$ and $r=1$. Using decomposition into partial fractions, we reduce to the case $f(x)=$ $(\mathrm{x}-\mathrm{b})^{-\mathrm{m}-1}, \quad|\mathrm{~b}|_{\mathrm{p}}>1$. But

$$
(\dot{x}-b)^{-m-1}=(-b)^{-m-1} \sum_{k=0}^{\infty}\binom{k+m}{m}\left(\frac{x}{b}\right)^{k},
$$

which converges on $D_{0}(1)$.
Lemma 4 (p-adic Cauchy integral formula). If $f$ is Krasner analytic in $D_{a}(r)$, and if $|\Gamma|_{p}=r$, then for fixed $z \in \Omega_{p}$

$$
\int_{a, \Gamma} \frac{f(x)(x-a)}{x-z} d x=\left\{\begin{array}{lll}
f(z) & \text { if } & |z-a|_{p}<r ;  \tag{3.2}\\
0 & \text { if } & |z-a|_{p}>r
\end{array}\right.
$$

In particular, this integral does not depend on the choice of $a$, $r$, or $r$ as long as $|z-a|_{p}$ remains either $<r$ or $>r$. More generally,

$$
\int_{a, \Gamma} \frac{f(x)(x-a)}{(x-z)^{m+1}} d x= \begin{cases}\frac{1}{m!} f^{(m)}(z) & \text { if }|z-a|_{p}<r  \tag{3.3}\\ 0 & \text { if }|z-a|_{p}>r\end{cases}
$$

Proof. By Leman 3 and the linearity and continuity of both sides (part (2) of Lemma 1), we reduce to the case $f(x)=(x-a)^{n}$.

$$
\begin{aligned}
& \text { Then write } \\
& \qquad \frac{1}{(x-z)^{m+1}}=\left\{\begin{array}{l}
\sum_{k=m+1}^{\infty}\binom{k-1}{m}(z-a)^{k-\mathbb{m}-1}(x-a)^{-k} \quad \text { if }|z-a|_{p}<r \\
(-1)^{m+1} \sum_{k=0}^{\infty}\binom{k+m}{m}(z-a)^{-k-m-1}(x-a)^{k} \quad \text { if }|z-a|_{p}>r .
\end{array}\right.
\end{aligned}
$$

Now use part (3) of Lemma 1 to conclude (3.3).
Lemma 5 (p-adic residue theorem). Let $f(x)=g(x) / h(x)$, where $g(x)$ is Krasner analytic in $D_{a}(r)$ (i.e., by Lemma 3, a power series) and $h(x)$ is a polynomial. Let $\left\{x_{i}\right\}$ be the roots of $h$ in $D_{a}(r)$, and suppose that all $\left|x_{i}-a\right|_{p}$ are strictly $<r$. Define res $x_{i} f$ to be the coefficient of $\left(x-x_{i}\right)^{-1}$ in the Laturent expansion of $f(x)$ at $x_{i}$. Then
$\int_{a, \Gamma} f(x)(x-a) d x=\sum \operatorname{res}_{x_{i}} f$.
Proof. Using the partial fraction decomposition of $1 / \mathrm{h}(\mathrm{x})$,
we reduce to the case $h(x)=\left(x-x_{i}\right)^{m+1}$. Then use (3.3) with $f(x)$ replaced by $g(x)$ and $z$ replaced by ' $x_{i}$.

The next lemma will be stated and proved in the form we shall need it, although some of the assumptions can be eliminated (the $D_{i}$ can have different radii, and $f(x)$ can approach a nonzero finite limit at infinity).

Lemma 6 (p-adic maximum modulus principle). Let $f(x)$ be a Krasner analytic function on $\Omega_{p}-U D_{i}$, where $D_{i}=D_{a_{i}}\left(r^{-}\right)$are open discs of radius $r$. Further suppose that $f(x) \longrightarrow 0$ as $|x|_{p}$ $\longrightarrow \infty$. Then $|f(x)|_{p}$ reaches its maximum on the boundary, i.e., if $|f(x)|_{p} \leq M$ for all $x$ with $\left|x-a_{i}\right|_{p}=r$ for some $i$, then $|f(x)|_{p} \leq M$ for all $x \in \Omega_{p}-U D_{i}$.

Proof. By the definition of Krasner analyticity, we immediateIy reduce to the case where $f(x)$ itself is a rational function with poles $b_{j} \in U_{D_{i}}$. Let $z \in \Omega_{p}$ be such that $\left|z-a_{i}\right|_{p}>x$ for all i. We must show that $|f(z)|_{p} \leq M$. Choose $r_{2}$ large enough so that $D_{0}\left(r_{2}\right)=D_{z}\left(r_{2}\right)$ contains $z$ and all of the $D_{i}$, and so that $|f(x)|_{p} \leq M$ for $|x|_{p}=r_{2}$. Let $\left|r_{2}\right|_{p}=r_{2}$. By part (1) of Lemma 1 ,

$$
\begin{equation*}
\left|\int_{z, \Gamma_{2}}^{1,} f(x) d x\right|_{p} \leq M \tag{3.4}
\end{equation*}
$$

On the other hand, by Lemma 5 ,

$$
\begin{align*}
\int_{z, \Gamma_{2}} f(x) d x & =\sum \operatorname{res} \frac{f(x)}{x-z} \\
& =f(z)+\sum_{j} \text { res }_{b_{j}} \frac{f(x)}{x-z} \tag{3.5}
\end{align*}
$$

Now let $|\Gamma|_{p}=r$. By Lemal 5, for each $i$

$$
\sum_{b_{j} \in D_{i}} \operatorname{res}_{b_{j}} \frac{f(x)}{x-z}=\int_{a_{i}, \Gamma} f(x) \frac{x-a_{i}}{x-z} d x .
$$

Since $|x-z|_{p}>\left|x-a_{i}\right|_{p}$ for $\left|x-a_{i}\right|_{p}=r$, it follows by part (1)
of Lemma 1 that for each $i$

$$
\left|\sum_{b_{j} \in D_{i}} \operatorname{res}_{b_{j}} \frac{f(x)}{x-z}\right|_{p} \leq \max _{\left|x-a_{i}\right|_{p}}=r|f(x)|_{p} \leq M .
$$

Combining this with (3.4) and (3.5) gives $|f(z)|_{p} \leq M$.
This concludes the basic lemas relating to the Shnirelman integral.

Let $\sigma \subset \Omega_{p}$ be a compact subset, such as $Z_{p}$ or $Z_{p}^{*}$. Let $\bar{\sigma}=\Omega_{p}-\sigma$ be its complement. For $z \in \bar{\sigma}$ let dist $(z, \sigma)$ denote the minimum of $|z-x|_{p}$ as $x$ ranges through $\sigma$.

Let $H_{0}(\bar{\sigma})$ denote the set of functions $\phi: \bar{\sigma} \longrightarrow \Omega_{p}$ which are Krasner analytic and vanish at infinity, i.e.,
(1) $\phi$ is a limit of rational functions whose poles are con-
tained in $\sigma$, the limit being uniform in any set $\bar{D}_{\sigma}(r)$ def $\left\{z \in \Omega_{p}\right\}$ $\operatorname{dist}(z, \sigma) \geq r\}$.

$$
\stackrel{(2)}{|z|_{\mathrm{p}}} \stackrel{\operatorname{iim}}{\rightarrow \infty} \phi(z)=0 .
$$

Remark. Strictly speaking, to say that $\phi$ is Krasner analytic on $\bar{\sigma}$ a priori means only that for every $r>0$ it is a uniform limit on $\bar{D}_{\sigma}(r)$ of rational functions $\phi_{n}$ with poles in $D_{\sigma}\left(r^{-}\right)$ dé $=\left\{z \in \Omega_{p} \mid\right.$ dist $\left.(z, \sigma)<r\right\}$. But if, for example, $\phi_{n}(z)=1 /(z-b)$ with $|b-a|_{p}=r_{1}<r$ for some $a \in \sigma$, then $\phi_{n}(z)=\sum_{j} \frac{(b-a)^{j}}{(z-a)^{j+1}}$ can be approximated uniformly on $\bar{D}_{\sigma}(r)$ by a rational function with pole $a \in \sigma$. Similarly if $\phi_{n}(z)=(z-b)^{-m} ;$ and any rational function $\phi_{\mathrm{n}}$ can be reduced to this case by partial fractions. Thus, the poles of $\phi_{\mathrm{n}}$ can be "shifted" to lie in $\sigma$.

Examples. 1. Since $\frac{d^{2}}{d z^{2}}\left((z+j)\left(\ln _{p}(z+j)-1\right)\right)=\frac{1}{z+j}$, we have

$$
\frac{d^{2}}{d z^{2}} G_{p}(z)=\lim _{n \rightarrow \infty} p^{-n} \sum_{0 \leq j<p} \frac{1}{z+j} \in H_{0}\left(\bar{z}_{p}\right),
$$

where $G_{p}$ is Diamond's $p$-adic $\log$ gamma function, see (7.4) in Chapter II. (It is not hard to show that the construction $\lim \mathrm{p}^{-\mathrm{n}} \sum f(\mathrm{z}+\mathrm{j})$ discussed in §II. 7 commutes with differentiation when $f$ is locally analytic.)
2. For any fixed $\xi \notin D_{1}\left(1^{-}\right)$, the first derivative of the twisted $\log$ gamma function $G_{p, \xi}$ (see (8.1) of Chapter II) is already Krasner analytic, since, by (8.6) of Chapter II,

$$
\frac{d}{d z} G_{p, \xi}(z)=-\int_{Z_{p}} \frac{d \mu_{\xi}(x)}{x+z}=-1 \lim \sum_{N \rightarrow \infty} \frac{1}{0 \leq a<p^{N}} \frac{\xi^{j}}{z+j} \frac{\xi^{p^{N}}}{1-H_{0}\left(\bar{Z}_{p}\right) .}
$$

3. If $\mu$ is any measure on $Z_{p}$ and $f(x) \in H_{0}\left(\bar{Z}_{p}\right)$, then it is not hard to check that $g(z)=\int_{Z_{p}} f(x-z) d \mu(x) \in H_{0}\left(\bar{Z}_{p}\right)$. In other words, $H_{0}\left(\bar{Z}_{p}\right)$ is stable under ${ }^{Z_{p}}$ convolution with measures on $Z_{p}$. The function $\frac{d}{d z} G_{p, \xi^{(-z)}}$ in Example 2 illustrates this.

For $r>0$ the set $D_{\sigma}\left(r^{-}\right)$is a finite (since $\sigma$ is compact) disjoint union of open discs of radius $r: D_{\sigma}\left(r^{-}\right)=U D_{a_{i}}\left(r^{-}\right)$. For example, if $\dot{\sigma}=z_{p}$ and $r=p^{-n}$ there are $p^{n+1}$ such discs.

Similarly, $D_{\sigma}(r) d \overline{\bar{E}} f\left\{z \in \Omega_{p}| | z-\left.a\right|_{p} \leq r\right.$ for sonie $\left.a \in \sigma\right\}$ is
a finite disjoint union of $D_{a_{i}}(r), a_{i} \in \sigma$.
Recall the definition $\bar{D}_{\sigma}(r)=\Omega_{p}-D_{\sigma}\left(r^{-}\right)=\left\{z \in \Omega_{p}\right\}|z-a|_{p} \geq r$ for all $a \in \sigma$ \}.

$$
\text { For } \phi \in H_{0}(\bar{\sigma}) \text { lèt }\|\phi\|_{r} \underset{\operatorname{def}}{=} \max _{z \in \bar{D}_{\sigma}(r)}^{\max }|\phi(z)|_{p} \text {. Obviously, }
$$

$\|\phi\|_{r_{1}} \geq\|\phi\|_{r}$ if $r_{1}<r$. By Lemma 6,

$$
\|\phi\|_{r}=\max _{\operatorname{dist}(z, \sigma)=r}|\phi(z)|_{p}
$$

We introduce a topology on $H_{0}(\bar{\sigma})$ by taking as a basis of open neighborhoods of zero
$U(r, \varepsilon)=\left\{\phi \mid\|\phi\|_{r}<\varepsilon\right\}$.
Note that $\|\phi\|_{r}$ is a continuous decreasing function of $r$. To see continuity, one easily reduces to the case when $\phi$ is a rational function with poles in $\sigma$, and then by partial fractions to the case when $\phi(x)=(x-a)^{-m}$, in which case it's obvious.

We next introduce a space of functions which will play a dual role to $H_{0}(\bar{\sigma})$ via a pairing defined using the Shnirelman integral.

Let
$\mathrm{B}_{\mathrm{r}}=\mathrm{B}_{\mathrm{r}}(\sigma)$ def $=\mathrm{ff}: \mathrm{D}_{\sigma}(\mathrm{r}) \longrightarrow \Omega_{\mathrm{p}} \mid \mathrm{f}$ is Krasner analytic on
By Lemma 3,

$$
\text { each } \left.D_{a_{i}}(r) \subset D_{\sigma}(r)\right\}
$$

$B_{r}=\left\{£ \mid\right.$ on each $D_{a_{i}}(r), \dot{f}$ is given by a convergent power series $\Sigma_{j} c_{i j}\left(z-a_{i}\right)^{j}$, i.e., $r^{j}\left|c_{i j}\right| \longrightarrow \longrightarrow \quad$ as
$j \longrightarrow \infty$ for each i\}.
If $r_{1}<r$, then restriction to $D_{\sigma}\left(r_{1}\right)$ gives an imbedding
$B_{r} \longrightarrow B_{r_{1}}$. We denote $L(\sigma)=\bigsqcup_{r>0} B_{r}$. By definition (see the beginning of $\$ I I .7), L(\sigma)$ is the set of locally analytic functions on $\sigma$.

$$
\begin{aligned}
& \text { For } f \in B_{r} \text { we set } \\
& \|f\|_{r ~ d e f} \overline{\max }_{i, j}\left|c_{i j}\right|_{p} r^{j}
\end{aligned}
$$

which is finite by definition. By Lemma 3, $\|f\|_{r}=\max _{z \in D_{\sigma}(r)}|f(z)|_{p}$.
Note that the inclusion $B_{r} \longrightarrow B_{r_{1}}$ for $r_{1}<r$ is continuous with respect to $\left\|\|_{r}\right.$ in $B_{r}$ and $\| \|_{r_{1}}$ in $B_{r_{1}}$.

Let $L^{*}(\sigma)$ be the set of continuous functionals on $L(\sigma)=$ $U_{B}{ }_{r}$, i.e., the set of linear maps $\mu$ (compatible with the restriction $\mathrm{B}_{\mathrm{r}} \longrightarrow \mathrm{B}_{\mathrm{r}_{1}}$ ) such that for all r
$\|\mu\|_{\mathrm{r} \text { déf }}=\underset{\substack{0 \neq f \in \mathrm{~B}_{\mathrm{r}}}}{ }|\mu(\mathrm{f})|_{\mathrm{P}} /\|\mathrm{f}\|_{\mathrm{r}}$
is finite. Note that $\|\mu\|_{r_{1}} \geq\|\mu\|_{r}$ if $r_{1}<r$. We do not require that $\|\mu\|_{r}$ remain bounded as $r \longrightarrow 0$.

Key example. Let $\mu$ be a measure on $\sigma$, i.e., a bounded additive map from compact-open subsets $U$ of $\sigma$ to $\Omega_{p}$. As in the case $\sigma=Z_{p}$ (see §II.2), the map

$$
\begin{equation*}
\mu: £ \longmapsto \int_{\sigma} f d \mu d=f \lim _{j} \sum_{i} f_{j}\left(U_{i j}\right) \mu\left(U_{i j}\right) \tag{3.6}
\end{equation*}
$$

(where $f_{j}$ is a sequence of locally constant functions which approach $f$ uniformly, and the $U_{i j}$ are compact-open sets on which $f_{j}$ is constant) is a well-defined functional on the continuous functions on $\sigma$, and a fortiori on $L(\sigma)$.

[^1]$\mu \in L^{*}(\sigma)$ has $\|\mu\|_{r} \leq M$ for all $r$. Define a function, also denoted $\mu$, on compact-open subsets $U$ of $\sigma$ by
$\mu(\mathrm{U})=\mu$ (characteristic function of U ).
(Note that any locally constant function is in $B_{r}$ for $r$ small.)
$\mu$ is obviously additive, and $|\mu(\mathrm{U})|_{\mathrm{p}} \leq\|\mu\|_{\mathrm{r}} \cdot \|$ char fn of $u \|_{\mathrm{r}}$
$=\|\mu\|_{r} \leq M$ for all $U$. This proves the lemma.
Choose P with $|\Gamma|_{p}=r$, and define
\[

$$
\begin{equation*}
f_{i j, r}(z)=\left(\frac{z-a_{i}}{F}\right)^{j} \text { restricted to } D_{a_{i}}(r) . \tag{3.7}
\end{equation*}
$$

\]

By Lemma 3, $B_{r}$ is the set of all series $f=\sum c_{i j} f_{i j, r}$ with
$c_{i j} \longrightarrow 0$ as $j \longrightarrow \infty$ for each of the (finitely many) $i$, and
$\|f\|_{r}=\max _{i, j}\left|c_{i j}\right|_{p}$. For $\mu \in L^{*}(\sigma)$ clearly

$$
\|\mu\|_{r}=\max _{i, j}\left|\mu\left(f_{i j, r}\right)\right|_{p}
$$

It can then be shown that the weak topology in $L^{*}(\sigma)$, which has basis of neighborhoods of zero

$$
\begin{equation*}
v_{f, \varepsilon \text { def }}=\left\{\left.\mu| | \mu(f)\right|_{p}<\varepsilon\right\} \tag{3.8}
\end{equation*}
$$

is equivalent to the (a priori stronger) topology having basis of neighborhoods of zero

$$
\begin{equation*}
V(r, \varepsilon) \underset{\text { def }}{\approx}\left\{\mu \mid\|\mu\|_{r}<\varepsilon\right\} . \tag{3.9}
\end{equation*}
$$

We shall prove this in the next section as a corollary of a general lemma on p-adic Banach spaces.

We shall of ten denote $\mu(f)$ by $(\mu(x), f(x))$.
Definition. For $\mu \in L^{*}(\sigma)$ the StieItjes transform $S \mu: \bar{\sigma} \longrightarrow$
$\Omega \mathrm{p}$ is the map
$\phi: z \longrightarrow\left(\mu(x), \frac{1}{z-x}\right)$.
We write $\phi=S \mu$.
Note that (3.10) makes sense, since for fixed $z \notin \sigma$, we have
$\frac{1}{z-x} \in B_{r}(\sigma)$ as soon as $r<\operatorname{dist}(z, \sigma)$.
Remark. If $\mu$ cones from a measure on $\sigma$ (also denoted $\mu$ ), then

$$
S \mu(z)=\int_{\sigma} \frac{d \mu(x)}{z-x} .
$$

This is slightly different from our earlier use of the term "Stie1tjes transform" for the dlog gamma type functions $\frac{d}{d z}$; namely, $\frac{d}{d z} G(z)=S i(-z)$.

Definition. For $\phi \in \mathrm{H}_{0}(\bar{\sigma})$ the Vishik transform $V \phi$ of $\phi$ is the functional on $\mathbf{L}(\sigma)=U B_{r}$ defined by

$$
\begin{equation*}
\mathrm{B}_{\mathrm{r}} \neq \mathrm{f} \longmapsto \sum_{i} \int_{a_{i}} \int_{\Gamma}^{L} \phi(x) f(x)\left(x-a_{i}\right) d x \tag{3.11}
\end{equation*}
$$

where this integral is the Shnirelman integral defined at the beginning of the section.

Lemma 8. (3.11) does not depend on the choice of centers $a_{i}$ or the choice of $\Gamma$ with $|\Gamma|_{p}=r$, and it is compatible with the inclusion $B_{r} \longrightarrow B_{r}$ for $r_{1}<r$.

Proof. For fixed $f$, the right side depends continuously on $\phi$ (with respect to $\|\phi\|_{r}$ ), so we may reduce to the case when $\phi$ is a rational function with poles in $\sigma$. In that case, by Lemma 5, the right side of (3.11) is simply $\Sigma$ res( $\phi f)$, and the lemma follows.

Remark. A function $\phi \in \mathrm{H}_{0}(\bar{\sigma})$ and a function $f \in \mathrm{~B}_{\mathrm{r}}(\sigma)$ have an annulus around $\sigma$ as a common domain of definition. The pairing ( $\phi, f$ ) $=V \phi(f)$ can be thought of as a pairing which evaluates the sum of the residues of the product. For example, if $\sigma$ is simply the point $\{0\}$, then $\phi(x)=\sum_{m<0} b_{m} x^{m}, f(x)=\sum_{n \geq 0} c_{n} x^{n}$, and
$(\phi, f)=$ coefficient of $\frac{1}{x}$ in $\phi(x) f(x)=\sum_{m+n=-1} b_{m} c_{n}$.
Theorem (Vishik), $V$ and $S$ are mutually inverse topological isomorphisms between $H_{0}(\bar{\sigma})$ and $L^{*}(\sigma)$. Under this isomorphism
the subspace $M(\sigma) \subset L^{*}(\sigma)$ of measures on $\sigma$ corresponds to the set of $\phi \in H_{0}(\bar{\sigma})$ such that $r\|\phi\|_{r}$ is bounded as $r \longrightarrow 0$.

Proof. Step 1. $S H \in H_{0}(\overline{\bar{\sigma}})$, and $S$ is continuous.
Notice that for fixed $z \in \bar{D}_{\sigma}(r)$ and for $r_{1}<r$ and $|\Gamma|_{p}=r_{1}$, the image of $\frac{1}{z-x}$ in $B_{r_{1}}$ is $\sum_{i, j}\left(z-a_{i}\right)^{-j-1} \Gamma^{j} f_{i j, r_{1}}$ ( $x$ ) (where $f_{i j, r_{1}}$ was defined in (3.7)). Then, since $\left|\left(\mu, f_{i j, r_{1}}\right)\right|_{p} \leq\|\mu\|_{r_{1}}$ $\left\|f_{i j, r_{1}}\right\|_{r_{1}}=\|\mu\|_{r_{1}}$, it follows that

$$
\left(\mu(x), \frac{1}{z-x}\right)=\lim _{n \rightarrow \infty} \sum_{i} \sum_{j<n} \frac{\Gamma^{j}}{\left(z-a_{i}\right)^{j+1}}\left(\mu, f_{i j}, r_{i}\right)
$$

is a uniform limit of rational functions with poles $a_{i} \in \sigma$ and value zero at infinity. At the same time we see continuity, since if $\mu \in V\left(r_{1}, \varepsilon\right)$, i.e., if $\|\mu\|_{r_{1}}<\varepsilon$, then
$\left\|\left(\mu(x), \frac{1}{z-x}\right)\right\|_{r}=\underset{z \in \bar{D}_{\sigma}(r)}{\max ^{(r)}}\left|\left(\mu(x), \frac{1}{z-x}\right)\right|_{p} \leq \underset{i, j}{\max } \frac{r_{1}{ }^{j}}{r^{j+1}}\|u\|_{r_{1}}$

$$
\begin{equation*}
<\frac{\varepsilon}{r} \tag{3.12}
\end{equation*}
$$

in other words, $S \mu \in U(r, \varepsilon / r)$.
Step 2. For $\phi \in H_{0}(\bar{\sigma})$, the functional $V \phi$ is continuous, i.e., $V$ maps $H_{0}(\bar{\sigma})$ to $L^{*}(\sigma)$.

Let $f \in \mathrm{~B}_{\mathbf{r}}(\sigma)$. Then

$$
\begin{aligned}
|(V \phi, f)|_{p} & =\left|\sum_{i} \int_{a_{i}, \Gamma} \phi(x) f(x)\left(x-a_{i}\right) d x\right|_{p} \\
& \leq \max _{i} \max _{\left|x-a_{i}\right|_{p}=r}\left|\phi(x) f(x)\left(x-a_{i}\right)\right|_{p}
\end{aligned}
$$

by part (1) of Lemma 1. But this is at most

$$
\max _{\operatorname{dist}(x, \sigma)=r}|\phi(x)|_{p} \max _{\operatorname{dist}(x, \sigma)=r}|f(x)|_{p}=r\|\phi\|_{r}\|f\|_{r} .
$$

Step 3. $V: H_{0}(\bar{\sigma}) \longrightarrow L *(\sigma)$ is continuous.
If $\phi \in U(r, \varepsilon)$, then we just saw that $|(V \phi, f)|_{P}<r \varepsilon\|f\|_{r}$
for $f \in B_{r}$. Thus $\|V \phi\|_{r}<r \varepsilon$, i.e, $V \phi \in V(r, r \varepsilon)$.
Step 4. VS = identity.
To see this, let $\mu \in L^{*}(\sigma), f \in B_{r}(\sigma)$, and denote $\phi=S \mu$. We must show that $(V \phi, f)=(\mu, f)$.

$$
\begin{aligned}
& \text { We have } \\
& \begin{aligned}
(V \phi, f) & =\sum_{i} \int_{a_{i}}\left(z-a_{i}\right) f(z) \phi(z) d z \\
& =\sum_{i} \int_{a_{i}}\left(z-a_{i}\right) f(z)\left(\mu(x), \frac{1}{z-x}\right) d z
\end{aligned}
\end{aligned}
$$

Since $\mu$ is linear and continuous, it commutes with the Shnirelman integral, and we have

$$
\begin{equation*}
(V \phi, f)=\left(\mu(x), \sum_{i} \int_{a_{i}, \Gamma} f(z) \frac{z-a_{i}}{z-x} d z\right) . \tag{3.13}
\end{equation*}
$$

Without loss of generality we may suppose that $r \notin\{|a-b| p \mid a, b$ $\epsilon \sigma\}$, in which case $r_{1}$ can be chosen less than $r$ but close enough so that $D_{\sigma}\left(r_{1}\right)$ is obtained from $D_{\sigma}(r)$ by merely shrinking each of the discs $D_{a_{i}}(r)$ (i.e., no elements of $\sigma$ are lost, so no new discs have to be added to cover $\sigma$ ). Now for $x \in D_{a_{i}}\left(r_{1}\right)$ the integral on the right in (3.13) is equal to $£(x)$, by Lemma
4. Thus, the restriction of $\sum_{i} \int_{a_{i}, r} f(z) \frac{z-a_{i}}{z-x} d z$ to $B_{r_{1}}$ is the same as the restriction of $f(x)$. Hence, $(V \phi, f)=(\mu, f)$.

Step 5. SV = identity.
Let $\phi \in H_{0}(\bar{\sigma})$.
We first suppose that $z$ is large, say $|z|_{p}>r=|\Gamma|_{p}$, where $r$ is taken large enough so that $\sigma \subset D_{0}(r)=D_{\sigma}(r)$. Let a be any point in $\sigma$. It is easy to see that $\phi(x)=\sum_{j=0}^{\infty} c_{j}(x-a)^{-j-1}$ for $x \in \bar{D}_{\sigma}(r)$ (as in the proof of Lemma 3).

Then

$$
\begin{aligned}
S V \phi(z) & =\left(V \phi(x), \frac{1}{z-x}\right)=\int_{a, \Gamma}(x-a) \phi(x) \frac{1}{z-x} d x \\
& =\int_{0, \Gamma} x \phi(x) \frac{1}{z-x} d \dot{x} \\
& =\int_{0,1 / \Gamma} \frac{1}{x} \phi\left(\frac{1}{x}\right) \frac{1}{z-1 / x} d x
\end{aligned}
$$

by the definition of the Shnirelman integral. Thus,

$$
S V \phi(z)=\frac{1}{z} \int_{0,1 / \Gamma} \phi\left(\frac{1}{x}\right) \frac{1}{x-1 / z} \mathrm{dx}
$$

$=\frac{1}{z} \mathrm{res}_{1 / \mathrm{z}} \frac{1}{\mathrm{x}} \frac{\phi(1 / \mathrm{x})}{\mathrm{x}-1 / \mathrm{z}}$ by Lemma 5
$=\phi(z)$.
(Alternately, we could expand $\phi(x)=\sum c_{j} x^{-j-1}$ and compute

$$
\begin{aligned}
S V \phi(z) & =\sum c_{j} \int_{0, \Gamma} x^{-j} \frac{1}{z-x} d x=\sum c_{j} \int_{0, \Gamma} \sum_{k=0}^{\infty} \frac{x^{k-j}}{z^{k+1}} d x \\
& \left.=\sum c_{j} / z^{j+1}=\phi(z) .\right)
\end{aligned}
$$

By Step 1, we know that $S V \phi(z) \in H_{0}(\bar{\sigma})$. Since $S V \phi(z)$ and $\phi(z)$ are both Krasner analytic on $\bar{\sigma}$ and agree on $\bar{D}_{0}(r)$, they nust agree everywhere.

$$
\begin{aligned}
& \text { Step 6. If } \mu \in \mathrm{E*}(\sigma) \text { and } \phi=S \mu \text {, then } \\
& \|\mu\|_{r}=r\|\phi\|_{r} .
\end{aligned}
$$

In Step 1 we saw that for any $r_{2}>r$
$\|\phi\|_{r_{2}} \leq \frac{1}{r_{2}}\|\mu\|_{r}$.
(We have replaced $r$ and $r_{1}$ by $r_{2}$ and $r$, respectively, in (3.12).) Letting $r_{2} \longrightarrow r$ and using the fact that $\|\phi\|_{r_{2}}$ is continuous in $r_{2}$, we obtain
$r\|\phi\|_{r} \leq\|\mu\|_{r}$.
On the other hand, in Step 2 we saw that $|(\mu, f)|_{p} \leq r\|\phi\|_{r}$. $\|f\|_{r}$ for all $f \in B_{r}$, and hence $\|\mu\|_{r} \leq r\|\phi\|_{r}$.

Step 7. $S(M(\sigma))=\left\{\phi \in H_{0}(\bar{\sigma}) \mid r\|\phi\|_{r}\right.$ is bounded $\}$.
This follows immediately from Step 6 and Lemma 7.
The proof of the theorem is now complete.
Remark. Amice and Vêlu [5] and Vishik [94] have studied socalled "h-admissible measures" on $\sigma$. These are elements $\mu \epsilon$ $L^{*}(\sigma)$ which instead of boundedness are required to satisfy the weaker condition

$$
r^{\mathrm{h}-j}\left|\mu\left(f_{i j, r}\right)\right|_{\mathrm{p}} \longrightarrow 0 \text { as } r \longrightarrow 0 \text { for all } i, j
$$

( $f_{i j, r}$ is the function (3.7)). For example, when $\sigma=z_{p}, j=0$, and $r=p^{-N}$, so that $f_{i j, r}$ is the characteristic function of some $a+{ }^{N_{Z}} Z_{p}$, this condition says that $\left|\mu\left(a+p_{Z_{p}}\right)\right|_{p}$ grows slower than $\mathrm{p}^{\mathrm{hN}}$. It is not hard to show that h-admissible measures $\mu$ correspond to functions $\phi \in H_{0}(\bar{\sigma})$ for which $r^{h+1}\|\phi\|_{\mathbf{r}}$. approaches zero as $\mathrm{r} \longrightarrow 0$.

Even the broader class of h-admissible measures are only a small part of $\mathrm{L} *(\sigma)$. For example, when $\sigma=\{0\}$, then $M(\sigma)$ is simply the constants, which correspond to elements $\phi \in H_{0}(\bar{\sigma})$ of the form $\phi(z)=\frac{\text { const }}{z}$. The h-admissible measures on $\{0\}$ correspond to the polynomials of degree at most $h$ in $1 / z$ (with no, constant term); while $L^{*}(\sigma)$ corresponds to all series $\sum c_{j} z^{-j}$ for which $r^{-j}\left|c_{j}\right|_{p} \longrightarrow 0$ as $j \longrightarrow \infty$ for every $r$.
4. p-adic spectral theorem

We start by discussing p-adic Banach spaces. For a more complete account, see [82]. In the process we fill in a technical gap in the last section, namely, we prove that in $L *(\sigma)$ the topology determined by

$$
\begin{equation*}
\nabla_{\mathrm{f}, \varepsilon}=\left\{\left.\mu| | \mu(\mathrm{f})\right|_{\mathrm{p}}<\varepsilon\right\} \tag{4.1}
\end{equation*}
$$

is equivalent to the topology determined by

$$
\begin{equation*}
\nabla(r, \varepsilon)=\left\{\mu \mid\|\mu\|_{r}<\varepsilon\right\} . \tag{4.2}
\end{equation*}
$$

Let $K$ be a field which is complete under a non-archimedean norm $\left|\left.\right|_{p}\right.$ (in'practice, $K$ will be a subfield of $\Omega_{p}$ ).

Definition. A Banach space over $K$ is a vector space $B$ supplied with a norm i| || from $B$ to the nonnegative real numbers such that for all $x, y \in B$ and $a \in K$; (1) $\|x\|=0$ if and only if $x=0$; (2) $\|x+y\| \leq \max (\|x\|,\|y\|)$; (3) $\|a x\|=|a|_{p}\|x\|$; (4) $B$ is complete with respect to \| \|.

We shall also assume that $\|B\|=|K| p$, i.e., for every $x \neq 0$ in $B$ there exists $a \in K$ such that $\|a x\|=1$.

By $\operatorname{Hom}\left(B_{1}, B_{2}\right)$ we mean the vector space of $K$-1inear continuous maps from $B_{1}$ to $B_{2} ; \operatorname{Hom}\left(B_{1}, B_{2}\right)$ is clearly a Banach space under the usual operator norm. We denote $\operatorname{End}(B)=\operatorname{Hom}(B, B)$.

If $B_{0}$ is a Banach space over $K_{0} \subset K$, by $B_{K}=B_{0} \hat{\otimes} R$ we mean the completed tensor product, i.e., the completion of the vector space $B_{0} \otimes_{K_{0}} K$.

Example. If $B=\left\{f=\Sigma c_{j} X^{j} \in Q_{p}[[x]]| | c_{j} \mid \longrightarrow 0\right\}$, with $\|f\|=\max _{j}\left|c_{j}\right|_{p}$, then $B_{\Omega_{p}}=\left\{f=\left.\sum c_{j} x^{j} \in \Omega_{p}[[x]]| | c_{j}\right|_{p} \longrightarrow 0\right\}$.

In practice, most interesting $\Omega_{p}$-Banach spaces $B$ are really defined over a finite extension $K$ of $Q_{p}$, in the sense that $B=B_{0} \hat{\otimes} \Omega_{p}$ for some $K$-Banach space $B_{0}$.

Canonical example. Let $J$ be any indexing set, and let $K(J)$ denote the set of all sequences $c=\left\{c_{j}\right\}_{j \in J}$ such that for every $\varepsilon$ only finitely many $\left|c_{j}\right|_{p}$ are $>\varepsilon$. Let $\|c\|=\max _{j}\left|c_{j}\right|_{p}$.

Note that $K(J) \hat{\otimes} \Omega_{p}=\Omega_{p}(\mathrm{~J})$.
Proposition. Let $K$ be a discrete valuation ring (for exam-
ple, a finite extension of $Q$, or the unramified closure of $Q_{p}$ ). Then any Banach space ${ }^{P}$ over $K$ is of the form (i.e., isomorphic to) $\mathrm{K}(\mathrm{J})$ for some J .

Proof. Let $0=0_{K}=\left\{\left.a \in K| | a\right|_{p} \leq 1\right\}, M=M_{K}=\left\{\left.a \in K| | a\right|_{p}<1\right\}$ $=\pi 0, \quad k=0 / M$. Let $E=\{x \in B \mid \quad\|x\| \leq 1\}, \bar{E}=E / \pi E$. Let $\left.\left\{e_{j}\right\}\right\}_{j \in J}$ be elements of $E$ whose reductions mod $\pi$ form a basis for the k-vector space $\bar{E}$. We claim that $B$ is isomorphic to $K(J)$. Given $x \in B$, find $a \in K$ such that $\|a x\| \leq 1$. Then for some $\left\{c_{1_{j}}\right\}_{j \in J}$ with $\left|c_{I_{j}}\right|_{p} \leq I$ and only finitely many $c_{1 j}$ nonzero we have: $a x-\sum c_{1 j} e_{j} \in \mathbb{T E}$. Repeating this process for $\frac{1}{\pi}$ (ax $\left.\sum c_{1 j} e_{j}\right)$, we successively find $a x=\sum_{j}\left(\sum_{i} \pi^{i} c_{i j}\right) e_{j}$. Let $c_{j}=$ $\frac{1}{a} \sum_{i} \pi^{i} c_{i j}$, and let $x$ correspond to $\left\{c_{j}\right\} \in K(J)$. Conversely, let every $\left\{c_{j}\right\} \in K(J)$ correspond to $\Sigma c_{j} e_{j}$. It is easy to see that this correspondence gives an isomorphism $B \simeq K(J)$.

Such a set $\left\{e_{j}\right\} \subset B$ is called a "Banach basis" for B.
Example. The space $\mathrm{B}_{\mathrm{r}}(\sigma)$ in the last section has Banach basis $f_{i j, r}$ (see (3.7)).

Corollary. If an $\Omega_{\mathrm{p}}$-Banach space $B$ is defined over a finite extension of $Q_{p}$, then $B$ is isomorphic to $\Omega_{p}(J)$ for some $J$.

Definition. The dual space $B^{*}$ of a K-Banach space $B$ is Hom( $B, K$ ), which is a Banach space with the usual operator norm.

Lemma 1. If $B \simeq K(J)$, then $B^{*}$ is isomorphic to the Banach
 In fact, if $\left\{e_{j}\right\}$ is a Banach basis, let $\left\{b_{j}\right\}$ be the map $\sum c_{j} e_{j} \longrightarrow \sum b_{j} c_{j}$. It is routine to check that this identifies $B *$ with the space of bounded sequences.

Definition. Let $B$ be a $K$-Banach space. A sequence $x_{1}, x_{2}$, $x_{3}, \ldots$ is said to be weakly convergent to $x$ if $h\left(x_{i}\right) \longrightarrow h(x)$ for all $h \in B^{*}$.

Lemma 2. Suppose $B \simeq K(J)$. If $x_{i} \longrightarrow x$ weakly, then
$\left\|x_{i}-x\right\| \longrightarrow 0$
Proof. Since any countable set of elements of $B$ is in the Banach subspace generated by a countable subset of our Banach basis for $B$, without loss of generality we may assume that $J$ is the positive integers. Replacing $x_{i}$ by $x_{i}-x$, without loss of generality we may suppose that $x=0$.

Suppose $\left\|x_{\mathbf{i}}\right\|$ does not approach zero. By passing to a subsequence, we may suppose that $\left\|x_{1}\right\|>\varepsilon$ for all $i$, We identify $B$ with $K(J)$, and let $x_{i}=\left\{a_{i j}\right\} \in K(J)$. Then $\left\|x_{i}\right\|=\max _{j}\left|a_{i j}\right|_{p}$ Let $\alpha_{i}$ denote the first $j$ for which $\left|a_{i j}\right|_{p}=\left\|x_{i}\right\|$, and let $\beta_{i}$ denote the last $j$ for which $\left|a_{i j}\right|_{p}=\left\|x_{i}\right\|$.

Case 1. $\alpha_{i}$ is bounded.
Then there exists some $j_{0}$ such that $\alpha_{i}=j_{0}$ for infinitely many i. Let $h \in B^{*}$ be the $j_{0}$-th coordinate map. Then for infinitely many $i$ we have

$$
\left|h\left(x_{i}\right)\right|_{p}=\left|a_{i j_{0}}\right|_{p}=\left|a_{i \alpha_{i}}\right|_{p}=\left\|x_{i}\right\|>\varepsilon,
$$

a contradiction.
Case 2. $\alpha_{i}$ is unbounded.
Let
$j_{0}=\alpha_{1}$
$j_{1}=\alpha_{i_{1}}$, where $i_{1}$ is chosen so that $\alpha_{i_{1}}>\beta_{1}$
$j_{2}=\alpha_{i_{2}}$, where $i_{2}$ is chosen so that $\alpha_{i_{2}}>\beta_{i_{1}}$
${ }^{2} \quad$
$j_{n}=\alpha_{i_{n}}$, where $i_{n}$ is chosen so that $\alpha_{i_{n}}>\beta_{i_{n-1}}$.
Let $h \in B^{*}$ be the sum of the $j_{n}$-th coordinate maps, i.e., $h\left(\left\{a_{j}\right\}\right)=\sum_{n} a_{j_{n}}$. Then for all $m$

$$
\left|h\left(x_{i_{m}}\right)\right|_{p}=\left|\sum_{\mathrm{n}}{a_{i_{m}} j_{n}}\right|_{p}=\left|a_{i_{m} i_{m}}\right|_{p}=\left\|x_{i_{m}}\right\| \mid
$$

and again $h\left(x_{i}\right)$ fails to approach 0 . This concludes the proof.

Corollary 1. If $B \simeq K(J)$ and $\left\{x_{i}\right\}, x_{i} \in B$, has the property that $h\left(x_{i}\right)$ approaches a finite limit for all $h \in B^{*}$, then $\left\{x_{i}\right\}$ converges in the norm to some x .

Proof. Let $y_{i}=x_{i}-x_{i+1}$. By Lemma 2, $\left\|y_{i}\right\| \| 0$. But then $\left\{x_{i}\right\}$ is a Cauchy sequence (since $\left\|x_{M}-x_{N}\right\| \leq \max _{M \leq i<N}\left\|x_{i}-x_{i+1}\right\|$ ), and the corollary follows from the completeness of B.

Corollary 2. The topologies on $L^{*}(\sigma)$ determined by (4.1) and (4.2) are equivalent.

Proof. Since the $V(x, \varepsilon)$-topology is trivially stronger than the $V_{f, \varepsilon}$-topology, it suffices to show that a sequence $\mu_{k}$ which converges to zero in the $V_{f}, \varepsilon^{\text {topology must converge to zero in }}$ the $V(r, \varepsilon)$-topology. Suppose that for every $f \in L(\sigma)=U B_{r}$ we. have $\mu_{k}(f) \longrightarrow 0$. We must show that for every $r,\left\|\mu_{k}\right\| \xrightarrow{r} 0$.

Without loss of generality we may suppose that $r \notin\left\{|a-b|_{p} \mid\right.$ $a, b \in \sigma\}$.

For any $r, \quad L^{*}(\sigma)$ maps to the dual $B_{r}^{*} \simeq \Omega_{p}(J)^{*}$, where $J$ indexes the $f_{i j, r}$. Namely, $\mu \longmapsto\left\{\mu\left(f_{i j, r}\right)\right\}_{i, j}$, and it is easy to see that the norm in ${ }_{B r}^{*}$ corresponds to $\|\mu\|_{r}$. Note that the image of $\mu$ has coordinates which approach zero as $j \longrightarrow \infty$ for each i. This is because, if we choose $r_{1}<r$ but large enough so that $\bigcup_{I} D_{a_{i}}\left(r_{1}\right)$ stili contains $\sigma$ (this can be done because $\sigma$ is compact, and no $b \in \sigma$ has $\left.\left|b-a_{i}\right|_{p}=r\right)$, then for all $i$, $j$

$$
\begin{aligned}
\left|\mu\left(f_{i j, r}\right)\right|_{p} & =\left|\mu\left(\left.f_{i j, r}\right|_{D_{\sigma}\left(r_{1}\right)}\right)\right|_{p}=\left|\mu\left(\left(\frac{\Gamma_{1}}{\Gamma}\right)^{j} f_{i j, r_{1}}\right)\right|_{p} \\
& =\left(\frac{r_{1}}{r}\right)^{j}\left|\mu\left(f_{i j, r_{1}}\right)\right|_{p} \leq\left(\frac{r_{1}}{r}\right)^{j}\|\mu\|_{r_{1}},
\end{aligned}
$$

which approaches zero as $j \longrightarrow \infty$. In other words, the compatibility requirement with $\mathrm{B}_{\mathrm{r}} \longrightarrow \mathrm{B}_{\mathrm{r}_{1}}$ forces $\mu$ to be a very special element of $\mathrm{B}_{\mathrm{r}}^{\mathrm{r}}$.

Now the subspace of $\mathrm{B}_{\mathrm{I}}^{*}$ consisting of elements whose coordin-
ates approach zero is of course isomorphic to $\Omega_{\mathrm{p}}(\mathrm{J})$. By Lemma 2, to show that $\left\|\mu_{\mathrm{k}}\right\| \longrightarrow 0$ it suffices to show that $\mathrm{g}\left(\mu_{\mathrm{k}}\right) \longrightarrow 0$ for all $g \in \Omega_{p}(J) *$. But if $g=\left\{g_{i j}\right\}$ under the isomorphism in Lemma 1, then
$g\left(\mu_{k}\right)=\sum_{i, j} g_{i j} \mu_{k}\left(f_{i j, r}\right)=\sum_{i, j} g_{i j}\left(\frac{\Gamma_{T}}{\Gamma}\right)^{j} \mu_{k}\left(f_{i j, r_{1}}\right)=\mu_{k}(f)$,
where $f=\sum g_{i j} \Gamma_{i}^{j} \Gamma^{-j} f_{i j, r_{1}} \in B_{r_{1}} . \quad$ And $\mu_{k}(f) \longrightarrow 0$ by assumption.

We now discuss operators $A \in \operatorname{End}(B)$.
If $B=K(J)$ with Banach basis $\left\{e_{j}\right\}_{j \in J}$, then $A$ corresponds to a matrix $\left\{a_{i j}\right\}_{i, j \in J}$ in the usual way:

$$
A e_{j}=\sum a_{i j} e_{i}
$$

It is easy to check that this is a norm-preserving isomorphism between End(B) and the Banach space of matrices $\left\{\mathrm{a}_{i j}\right\}$ having finite $\left\|\left\{a_{i j}\right\}\right\| \|=\max _{i, j}\left|a_{i j}\right|_{p}$ and having the property that for each $j, a_{i} \longrightarrow 0$ as $i \longrightarrow \infty$. In other words, when $B=K(J)$, A can be thought of as a matrix whose columns are in $B$ and whose rows are in $\mathrm{B}^{*}$.

An operator A is said to be "completely continuous" if it can be approximated by operators having finite-dimensional image. In terms of matrices, this means that $a_{i j} \longrightarrow 0$ as $i \longrightarrow \infty$ uniformily in $j$; in other words, the norm of the i-th row approaches zero. Such operators occur in Dwork's theory (e.g. [25]), and in [82] Serre gives a Riesz and Fredholm theory for them.

However, many simple operators are not completely continuous: the identity operator, for example, or the operator $\left(\frac{x}{d x}\right)^{n}$ on $\left\{\Sigma c_{i} x^{i} \mid c_{i} \longrightarrow 0\right\}$ (which has diagonal matrix $a_{i i}=i^{n}$ ).

For simplicity, we shall assume our Banach spaces are of the
form $\Omega_{p}(J)$. As mentioned before, all $\Omega_{p}$-Banach spaces which are defined over a discrete valuation subfield of $\Omega_{p}$ are of this form.

Definition. For $A \in \operatorname{End}(B)$ let $\sigma_{A}=\left\{\lambda \in \Omega_{p} \mid A-\lambda\right.$ does not have an inverse in $\operatorname{End}(B)\}$ denote the spectrum of $A$.

Definition. An operator $A \in \operatorname{End}(B)$ is called analytic if the "resolvent" operator $R_{A}(z)=(z-A)^{-1}$ is Krasner analytic in the complement of $\sigma_{A}$, in the sense that for all $x \in B$ and $h \in B *$, $h\left(R_{A}(z) x\right)$ as a function of $z$ lies in $H_{0}\left(\bar{\sigma}_{A}\right)$. If $B=\Omega_{p}(J)$, then in terms of matrices this is equivalent to the condition that each matrix entry in $R_{A}(z)$ be a Krasner analytic function of $z$ on $\bar{\sigma}_{A}$ (and vanish at infinity).

Vishik's spectral theory applies to analytic operators A whose spectrum $\sigma_{A}$ is a compact subset of $\Omega_{\mathrm{P}}$.

Example. $\mathrm{x} \frac{\mathrm{d}}{\mathrm{dx}}$ acting on $\mathrm{B}=\left\{\sum \mathrm{c}_{\mathrm{i}} \mathrm{x}^{\mathbf{i}} \mid \mathrm{c}_{\mathrm{i}} \longrightarrow 0\right\}$ has spectrum $Z_{p}$, and its resolvent is Krasner analytic on $\bar{Z}_{p}$.

It is possible for A to have a compact spectrum but not satisfy the analyticity condition. Here is an example of Vishik where the spectrum is empty. (Since the only Krasner analytic functions on all of $\Omega_{p}$, by Lemma 3 of 53 , are everywhere convergent power series, and since only the zero power series has value 0 at infinity, it follows that in the case of an empty spectrum $R_{A}(z)$ has no chance of having matrix entries in $H_{0}\left(\bar{\sigma}_{A}\right)$.)

Example. Let $B$ be the set of $\left\{a_{i}\right\}_{i \in Z}$ such that $\left\|\left\{a_{i}\right\}\right\|_{d e f}$ $\max _{i \in Z}\left|a_{i}\right|_{p}$ is finite and $a \longrightarrow 0$ as $i \longrightarrow-\infty$. Let $A$ be the shift operator $A\left(\left\{a_{i}\right\}\right)=\left\{b_{i}\right\}$ where $b_{i}=a_{i+1}$.

Claim. For all $z \in \Omega_{p},(z-A)$ has a continuous inverse $f_{z}$.
Proof of claim. We want to find $f_{z}:\left\{b_{j}\right\} \longmapsto\left\{a_{i}\right\}$. such that
$a_{i}=f_{z, i}\left(\left\{b_{j}\right\}\right)$ satisfies $\quad z a_{i}-a_{i+1}=b_{i}$.
Case (1). $|z|_{p} \leq 1$. Set $a_{i}=-b_{i-1}-z b_{i-2}-z^{2} b_{i-3}-\cdots$,
which converges because $b \longrightarrow 0$ as $j \longrightarrow-\infty$. Clearly $\left\{a_{i}\right\} \in B ;$ $z a_{i}-a_{i+1}=b_{i}$; and this map $f_{z}$ is continuous.

Case (2). $|z|_{p}>1$. Set,$a_{i}=b_{i} z^{-1}+b_{i+1} z^{-2}+\cdots$, Again $\left\{a_{i}\right\} \in B ; \quad z a_{i}-\dot{a}_{i+1}=b_{i} ;$ and the map is continuous.

Let. B be a Banach space of the form $\Omega_{\mathrm{p}}(J)$. Let $F(\mathrm{x})$ be an analytic operator-valued function on the complement of a compact set $\sigma$, i.e., for all $y \in B$ and $h \in B^{*}$ the $\Omega_{p}$-valued function $\mathrm{F}_{\mathrm{h}, \mathrm{y}}(\mathrm{x}) \quad \mathrm{d} \overline{\mathrm{e}}_{\mathrm{f}} \mathrm{h}(\mathrm{F}(\mathrm{x}) \mathrm{y})$
belongs to $H_{0}(\bar{\sigma})$. Let $a \in \sigma,|\Gamma|=r$, and suppose that there are no $b \in \sigma$ such that $|b-a|_{p}=r$.

Definition, Let $S_{n}=\frac{1}{n} \sum_{\xi^{n}=1} F(a+\xi!)$. Then

Lemma 3. The limit (4.3) converges in the operator norm to a continuous operator.

Proof. Let $y \in B$. For all $h \in B^{*}$, since $F_{h, y} \in H_{0}(\bar{\sigma})$, it follows that the ordinary Shnirelman integral $\int_{a, \Gamma} F_{h, y}(x) d x$ exists. That is, $h\left(S_{n} y\right)$ approaches a finite limit for all $h$. By Corollary 1 to Lemma 2, $S_{n} y$ converges in the norm. By the uniform boundedness principle, $S_{n}$ converges to a continuous operator.

Note that from the proof of Lemma 3 it follows that

$$
\begin{equation*}
h\left(\int_{a, \Gamma} F(x) d x y\right)=\int_{a, \Gamma} F_{h, y}(x) d x . \tag{4.4}
\end{equation*}
$$

Spectral theorem (Vishik). Let $B \simeq \Omega_{p}(J)$, and let $A \in \operatorname{End}(B)$ be analytic with compact spectrum $\sigma_{A}$. Then the operator-valued distribution
$\mu_{A} \quad \underset{\operatorname{def}}{ } \quad V R_{A}$
where $V$ is the Vishik transform in $\S 3$, gives a continuous homomorphism from the algebra $L\left(\sigma_{A}\right)$ to the algebra End $(B)$. For $f \in$ ${ }^{B}{ }_{r}\left(\sigma_{A}\right)$, the operator $\mu_{A}(£)$ is defined as

$$
\sum_{i} \int_{a_{i}, r} f(x)\left(x-a_{i}\right)(x-A)^{-1} d x
$$

where $D_{\sigma}(r)=U D_{a_{i}}(r)$ is a covering of $\sigma$ by discs of radius $r$. In addition, the following inversion formula holds:

$$
\begin{equation*}
R_{A}(z)=\left(\mu_{A}(x), \frac{1}{z-x}\right) \tag{4.5}
\end{equation*}
$$

Corollary 1. For all $j \geq 0, A^{j}=\left(\mu_{A}(x), x^{j}\right)$.
Proof of corotlary. For any fixed $z$ with $|z|_{p}>\max _{x \in 0}|x|_{p}$ and $|z|_{p}>\|A\|$, we can write $\frac{1}{z-x}=\sum_{j=0}^{\infty} z^{-j-1} x^{j} \ln (4.5)$. Since $|z|_{P}>\|A\|$, we also have $R_{A}(z)=\sum_{j=0}^{\infty} z^{-j-1} A^{j}$. By the continuity of $\mu_{A}: L\left(\sigma_{A}\right) \longrightarrow$ End $(B)$, this gives us $\sum z^{-j-1} A^{j}=$ $\sum z^{-j-1}\left(\mu_{A}(x), x^{j}\right)$. Since this holds for all large $z$, the coefficients can be equated, and the corollary is proved.

Remark. For $j=1$, if we write ( $\left.\mu_{A}, f\right)$ using the $\int$ notation, we obtain the usual form for a spectral theorem:
$A=\int_{\sigma_{A}} x d \mu_{A}(x)$.
Proof of the spectral theorem. First of all, it is easy to see that $\mu_{A}(f)$ is a bounded operator, and that $\mu_{A}: I\left(\sigma_{A}\right) \longrightarrow \operatorname{End}(B)$ is continuous. The key assertion is multiplicativity:

$$
\mu_{A}\left(f_{1} f_{2}\right)=\mu_{A}\left(f_{1}\right) \mu_{A}\left(f_{2}\right) \quad \text { for } \quad f_{1}, f_{2} \in L\left\langle\sigma_{A}\right)
$$

Suppose that $f_{1}, f_{2} \in B_{r}$. Let $r>r^{\prime}>r_{1}, \quad\left|\Gamma^{\dagger}\right|_{p}=r^{\prime},\left|r_{1}\right|_{p}$ $=r_{1}$. We can choose $r^{\prime}$ so that $r^{\prime} \notin\left\{|a-b|_{p} \mid a, b \in \sigma_{A}\right\}$, in which case $r_{1}$ can be chosen so that $D_{\sigma_{A}}\left(r_{1}\right)$ is obtained from $D_{\sigma_{A}}\left(r^{\prime}\right)=U_{D_{a_{i}}}\left(r^{\prime}\right)$ by merely shrinking each disc. Thus, $D_{\sigma_{A}}\left(r_{1}\right)$.
$=U D_{a_{i}}\left(r_{1}\right) . \quad$ Now

$$
\mu_{A}\left(f_{1}\right) \mu_{A}\left(f_{2}\right)=\sum_{i} \sum_{j} \int_{a_{j}, \Gamma^{\prime}} \int_{a_{i}, r_{1}} \frac{\left(x^{\prime}-a_{j}\right)\left(x-a_{i}\right) f_{2}\left(x^{\prime}\right) f_{1}(x)}{\left(x^{\prime}-A\right)(x-A)} d x d x^{\prime}
$$

Since $\frac{1}{\left(x^{\prime}-A\right)(x-A)}=\frac{1}{x-x^{\prime}}\left(\left(x^{\prime}-A\right)^{-1}-(x-A)^{-1}\right)$, this equals

$$
\begin{aligned}
& \sum_{i} \sum_{j} \int_{a_{j}, \Gamma^{\prime}}\left(x^{\prime}-a_{j}\right) f_{2}\left(x^{\prime}\right)\left(x^{\prime}-A\right)^{-1} \int_{a_{i}, \Gamma_{1}} \frac{\left(x-a_{i}\right) f_{1}(x)}{x-x^{\prime}} d x d x^{\prime} \\
& +\sum_{i} \sum_{j} \int_{a_{i}, r_{1}}\left(x-a_{i}\right) f_{1}(x)(x-A)^{-1} \int_{a_{j}, \Gamma^{\prime}} \frac{\left(x^{\prime}-a_{j}\right) f_{2}\left(x^{\prime}\right)}{x^{\prime}-x} d x^{\prime} d x .
\end{aligned}
$$

But by Lemma 4 of the last section, the inner integral in the first sum is zero, and the inner integral in the second sum is zero for $j \neq i$ and is $f_{2}(x)$ for $j=i$. Thus,

$$
\mu_{A}\left(f_{1}\right) \mu_{A}\left(f_{2}\right)=\sum_{i} \int_{a_{i}, \Gamma_{1}}\left(x-a_{i}\right) f_{1}(x) f_{2}(x)(x-A)^{-1} d x=\mu_{A}\left(f_{1} f_{2}\right)
$$

Finally, to prove the inversion formula, for any $y \in B$ and $h \in B^{*}$ we have

$$
h\left(R_{A}(z) y\right)=R_{A, h, y}(z)=S V R_{A, h, y}(z) \quad \text { by the theorem in } £ 3
$$

$$
=\left(V_{R_{A}, h, y}(x), \frac{1}{z-x}\right)^{=}=h\left(\left(\mu_{A}(x), \frac{1}{z-x}\right) y\right)
$$

Thus, (4.5) is an immediate consequence of the theorem on inverting the p-adic Stieltjes transform. This completes the proof of the theorem.

Corollary 2. Under the conditions of the theorem, the following two conditions are equivalent: (1) dist $(z, \sigma)\left\|R_{A}(z)\right\|$ is bounded as $z \longrightarrow \sigma$; (2) $\mu_{A}$ is a projection-valued measure, i.e., a bounded homomorphism from the Boolean algebra of compactopen subsets of $\sigma_{A}$ to the algebra End (B). In this case
$\max _{z \in \bar{\sigma}_{A}}\left(\operatorname{dist}(z, \sigma)\left\|R_{A}(z)\right\|\right)=\max _{U}\left\|\mu_{A}(U)\right\|$.
The proof is exactly like Steps 6 and 7 in the proof of the theorem in $\S 3$.

Corollary 3. Let $B$ and $A \in E \operatorname{Ad}(B)$ be as in the theorem. For a11 $r>0$ the resolvent $R_{A}(z)$ can be uniformly approximated on $\bar{D}_{\sigma_{A}}(r)$ by rational functions $\sum_{i} \sum_{j=1}^{N} A_{i j}\left(z-a_{i}\right)^{-j}$ with operator coefficients and with poles in $\sigma_{A}$.

The proof is just like Step 1 in the proof of the theorem in $\S 3$.

Remarks. 1. One could alternately take Corollary 3 as the definition of an analytic operator, in which case the spectral theorem would hold for an arbitrary Banach space. However, in practice the "weaker" definition is often easier to check than the strong condition in the corollary.
2. The operators in Corollary 2 are the closest p-adic analogs of normal operators or operators of scalar type [24] in a Hilbert space.
3. It is not hard to prove that operators for which
$\operatorname{dist}\left(z, \sigma_{A}\right)^{h+1}\left\|R_{A}(z)\right\| 0$ as $\quad z \longrightarrow \sigma$ correspond to "h-admissible" $\mu_{A}$ (see the remark at the end of the last section).
4. Vishik has also proved a generalization to functions of (the spectra of) several analytic operators. Namely, let $A_{1}, \ldots$, $A_{n} \in \operatorname{End}(B)$ be commuting analytic operators with compact spectra $\sigma_{A_{1}}, \ldots, \sigma_{A_{n}}$. Let $\sigma=\sigma_{A_{1}} \times \cdots{ }_{A_{n}} \subset \Omega_{p}^{n}$, and use the completed tensor product to define $B_{r}(\sigma)$ and $L(\sigma): B_{r}(\sigma)=B_{r}\left(\sigma_{A_{1}}\right) \hat{\otimes} \cdots$
 be the continuous homomorphism from the algebra $\mathrm{L}(\sigma)$ to the algebra End(B) which is made up from the $\mu_{A_{i}}$ in the theorem. For $f \in L(\sigma)$ denote $f\left(A_{1}, \ldots, A_{n}\right)=\mu(f)$. Then Vishik shows that $f\left(A_{1}, \ldots, A_{n}\right)$ is an analytic operator with compact spectrum

$$
\begin{aligned}
& \sigma_{f\left(A_{1}, \ldots, A_{n}\right)} \subset f\left(\sigma_{A_{1}}, \ldots, \sigma_{A_{n}}\right), \text { and for } z \notin f\left(\sigma_{A_{1}}, \ldots, \sigma_{A_{n}}\right) \text { we } \\
& \text { have } \\
& \quad R_{f\left(A_{1}, \ldots, A_{n}\right)}(z)=\left(\mu\left(x_{1}, \ldots, x_{n}\right), \frac{1}{z-f\left(x_{1}, \ldots, x_{n}\right)}\right) .
\end{aligned}
$$

$$
\text { In addition, } \mu_{f\left(A_{1}, \ldots, A_{n}\right)}=f_{*} \mu \text {, i.e., if } \ell \in L\left(f\left(\sigma_{A_{1}}, \ldots, \sigma_{A_{n}}\right)\right) \text {, }
$$

then

$$
\mu_{f\left(A_{1}, \ldots, A_{n}\right)^{(l)}=\mu(\ell \circ f) .}
$$

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[^0]:    Proposition. For K totally real, Leopoldt's confecture is

[^1]:    Lemma 7. $\mu \in \mathrm{L}^{*}(\sigma)$ comes from a measure on $\sigma$ if and only if $\|\mu\|_{r} \xrightarrow{\text { is bounded as }} r \longrightarrow 0$.

    Proof. Using (3.6), it is easy to check that $\underset{r \rightarrow 0}{\lim }\|\mu\|_{r}=$ $\max _{U}|\mu(U)|_{P}<\infty$ whenever $\mu$ is a measure. Conversely, suppose

