# Arithmeticity in the Theory of Automorphic Forms 

Goro Shimura

## Selected Titles in This Series

82 Goro Shimura, Arithmeticity in the theory of automorphic forms, 2000
81 Michael E. Taylor, Tools for PDE: Pseudodifferential operators, paradifferential operators, and layer potentials, 2000
80 Lindsay N. Childs, Taming wild extensions: Hopf algebras and local Galois module theory, 2000
79 Joseph A. Cima and William T. Ross, The backward shift on the Hardy space, 2000
78 Boris A. Kupershmidt, KP or mKP: Noncommutative mathematics of Lagrangian, Hamiltonian, and integrable systems, 2000
77 Fumio Hiai and Dénes Petz, The semicircle law, free random variables and entropy, 2000
76 Frederick P. Gardiner and Nikola Lakic, Quasiconformal Teichmüller theory, 2000
75 Greg Hjorth, Classification and orbit equivalence relations, 2000
74 Daniel W. Stroock, An introduction to the analysis of paths on a Riemannian manifold, 2000
73 John Locker, Spectral theory of non-self-adjoint two-point differential operators, 2000
72 Gerald Teschl, Jacobi operators and completely integrable nonlinear lattices, 1999
71 Lajos Pukánszky, Characters of connected Lie groups, 1999
70 Carmen Chicone and Yuri Latushkin, Evolution semigroups in dynamical systems and differential equations, 1999
69 C. T. C. Wall (A. A. Ranicki, Editor), Surgery on compact manifolds, second edition, 1999
68 David A. Cox and Sheldon Katz, Mirror symmetry and algebraic geometry, 1999
67 A. Borel and N. Wallach, Continuous cohomology, discrete subgroups, and representations of reductive groups, second edition, 2000
66 Yu. Ilyashenko and Weigu Li, Nonlocal bifurcations, 1999
65 Carl Faith, Rings and things and a fine array of twentieth century associative algebra, 1999
64 Rene A. Carmona and Boris Rozovskii, Editors, Stochastic partial differential equations: Six perspectives, 1999
63 Mark Hovey, Model categories, 1999
62 Vladimir I. Bogachev, Gaussian measures, 1998
61 W. Norrie Everitt and Lawrence Markus, Boundary value problems and symplectic algebra for ordinary differential and quasi-differential operators, 1999
60 Iain Raeburn and Dana P. Williams, Morita equivalence and continuous-trace $C^{*}$-algebras, 1998
59 Paul Howard and Jean E. Rubin, Consequences of the axiom of choice, 1998
58 Pavel I. Etingof, Igor B. Frenkel, and Alexander A. Kirillov, Jr., Lectures on representation theory and Knizhnik-Zamolodchikov equations, 1998
57 Marc Levine, Mixed motives, 1998
56 Leonid I. Korogodski and Yan S. Soibelman, Algebras of functions on quantum groups: Part I, 1998
55 J. Scott Carter and Masahico Saito, Knotted surfaces and their diagrams, 1998
54 Casper Goffman, Togo Nishiura, and Daniel Waterman, Homeomorphisms in analysis, 1997
53 Andreas Kriegl and Peter W. Michor, The convenient setting of global analysis, 1997

For a complete list of titles in this series, visit the AMS Bookstore at www.ams.org/bookstore/.

This page intentionally left blank

# Arithmeticity in the Theory of Automorphic Forms 

## Goro Shimura

Editorial Board<br>Georgia Benkart<br>Peter Landweber<br>Michael Loss<br>Tudor Ratiu, Chair

2000 Mathematics Subject Classification. Primary 11Fxx, 14K22, 14K25, 32Nxx, 32A99, 32W99.

Abstract. The main theme of the book is the arithmeticity of the critical values of certain zeta functions associated with algebraic groups. The author also included some basic material about arithmeticity of modular forms in general and a treatment of analytical properties of zeta functions and Eisenstein series on symplectic groups. The book can be viewed as a companion to the previous book, Euler Products and Eisenstein Series, by the same author (AMS, 1997).

For researchers and graduate students working in number theory and the theory of modular forms.

## Library of Congress Cataloging-in-Publication Data

Shimura, Goro, 1930-
Arithmeticity in the theory of automorphic forms / Goro Shimura.
p. cm. - (Mathematical surveys and monographs, ISSN 0076-5376; v. 82)

Includes bibliographical references and index.
ISBN 0-8218-2671-9 (alk. paper)

1. Automorphic forms. I. Title. II. Mathematical surveys and monographs ; no. 82.

## QA331.S4645 2000

512'.7-dc21
00-032273

AMS softcover ISBN: 978-0-8218-4961-3

Copying and reprinting. Individual readers of this publication, and nonprofit libraries acting for them, are permitted to make fair use of the material, such as to copy a chapter for use in teaching or research. Permission is granted to quote brief passages from this publication in reviews, provided the customary acknowledgment of the source is given.

Republication, systematic copying, or multiple reproduction of any material in this publication is permitted only under license from the American Mathematical Society. Requests for such permission should be addressed to the Acquisitions Department, American Mathematical Society, 201 Charles Street, Providence, Rhode Island 02904-2294 USA. Requests can also be made by e-mail to reprint-permission@ams.org.
(C) 2000 by the American Mathematical Society. All rights reserved.

Reprinted by the American Mathematical Society, 2010.
The American Mathematical Society retains all rights except those granted to the United States Government.

Printed in the United States of America.
(a) The paper used in this book is acid-free and falls within the guidelines
established to ensure permanence and durability.
Visit the AMS home page at http://www.ams.org/

## TABLE OF CONTENTS

Preface ..... vii
Notation and Terminology ..... ix
Introduction ..... 1
Chapter I. Automorphic Forms and Families of Abelian Varieties ..... 7

1. Algebraic preliminaries ..... 7
2. Polarized abelian varieties ..... 13
3. Symmetric spaces and factors of automorphy ..... 17
4. Families of polarized abelian varieties ..... 22
5. Definition of automorphic forms ..... 30
6. Parametrization by theta functions ..... 37
Chapter II. Arithmeticity of Automorphic Forms ..... 45
7. The field $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ ..... 45
8. Action of certain elements of $\widetilde{G}_{\mathbf{A}}$ on $\mathfrak{K}$ ..... 50
9. The reciprocity-law at CM-points and rationality of auto- morphic forms ..... 58
10. Automorphisms of the spaces of automorphic forms ..... 67
11. Arithmeticity at CM-points ..... 76
Chapter III. Arithmetic of Differential Operators and Nearly Holomorphic Functions ..... 87
12. Differential operators on symmetric spaces ..... 87
13. Nearly holomorphic functions ..... 98
14. Arithmeticity of nearly holomorphic functions ..... 107
15. Holomorphic projection ..... 122
Chapter IV. Eisenstein Series of Simpler Types ..... 127
16. Eisenstein series on $U\left(\eta_{n}\right)$ ..... 127
17. Arithmeticity and near holomorphy of Eisenstein series ..... 137
18. Eisenstein series in the Hilbert modular case ..... 149
Chapter V. Zeta Functions Associated with Hecke Eigenforms ..... 159
19. Formal Euler products and generalized Möbius functions ..... 159
20. Dirichlet series obtained from Hecke eigenvalues and Fourier coefficients ..... 166
21. The Euler products for the forms of half-integral weight ..... 175
22. The largest possible pole of $\mathcal{Z}(s, \mathbf{f}, \chi)$ ..... 177
Chapter VI. Analytic Continuation and Near Holomor- phy of Eisenstein Series of General Types ..... 185
23. Eisenstein series of general types ..... 185
24. Pullback of Eisenstein series ..... 194
25. Proof of Theorems in Sections 20 and 23 ..... 203
26. Near holomorphy of Eisenstein series in Case UB ..... 208
Chapter VII. Arithmeticity of the Critical Values of Zeta Functions and Eisenstein Series of General Types ..... 219
27. The spaces of holomorphic Eisenstein series ..... 219
28. Main theorems on arithmeticity in Cases SP and UT ..... 230
29. Main theorems on arithmeticity in Case UB ..... 240
Appendix ..... 247
A1. The series associated to a symmetric matrix and Gauss sums ..... 247
A2. Metaplectic groups and factors of automorphy ..... 251
A3. Transformation formulas of general theta series ..... 262
A4. The constant term of a theta series at each cusp depends only on the genus ..... 272
A5. Theta series of a hermitian form ..... 274
A6. Estimate of the Fourier coefficients of a modular form ..... 278
A7. The Mellin transforms of Hilbert modular forms ..... 282
A8. Certain unitarizable representation spaces ..... 285
References ..... 297
Index ..... 301

## PREFACE

A preliminary idea of writing the present book was formed when I gave the Frank J. Hahn lectures at Yale University in March, 1992. The title of the lectures was "Differential operators, nearly holomorphic functions, and arithmetic." By "arithmetic" I meant the arithmeticity of the critical values of certain zeta functions, and I talked on the results I had on $G L_{2}$ and $G L_{2} \times G L_{2}$. At that time the American Mathematical Society wrote me that they were interested in publishing my lectures in book form, but I thought that it would be desirable to discuss similar problems for symplectic groups of higher degree. Though I had satisfactory theories of differential operators and nearly holomorphic functions applicable to higher-dimensional cases, our knowledge of zeta functions on such groups was fragmentary and, at any rate, was not sufficient for discussing their critical values. Therefore I spent the next few years developing a reasonably complete theory, or rather, a theory adequate enough for stating general results of arithmeticity that cover the cases of all congruence subgroups of a symplectic group over an arbitrary totally real number field, including the case of half-integral weight.

On the other hand, I had been interested in arithmeticity problems on unitary groups for many years, and in fact had investigated some Eisenstein series on them. Therefore I thought that a book including the unitary case would be more appealing, and I took up that case as a principal topic of my NSF-CBMS lectures at the Texas Christian University in May, 1996. The expanded version of the lectures was eventually published by the AMS as "Euler products and Eisenstein series."

After this work, I felt that the time was ripe for bringing the original idea to fruition, which I am now attempting to do in this volume. To a large extent the present book may be viewed as a companion to the previous one just mentioned, and our arithmeticity concerns that of the Euler products and Eisenstein series discussed in it; I did not include the cases of $G L_{2}$ and $G L_{2} \times G L_{2}$. Those cases are relatively well understood, and it is my wish to present something new. Though the arithmeticity in that sense is the main new feature, as will be explained in detail in the Introduction, I have also included some basic material concerning arithmeticity of modular forms in general, and also a treatment of analytic properties of zeta functions and Eisenstein series on symplectic groups which were not discussed in the previous book.

It is a pleasure for me to express my thanks to Haruzo Hida, who read the manuscript and contributed many useful suggestions.

Princeton
February, 2000
Goro Shimura

This page intentionally left blank

## NOTATION AND TERMINOLOGY

We denote by $\mathbf{Z}, \mathbf{Q}, \mathbf{R}$, and $\mathbf{C}$ the ring of rational integers, the fields of rational numbers, real numbers, and complex numbers, respectively. We put

$$
\mathbf{T}=\{z \in \mathbf{C}| | z \mid=1\}
$$

We denote by $\overline{\mathbf{Q}}$ the algebraic closure of $\mathbf{Q}$ in $\mathbf{C}$, and for an algebraic number field $K$ we denote by $K_{\mathrm{ab}}$ the maximum abelian extension of $K$. If $p$ is a rational prime, $\mathbf{Z}_{p}$ and $\mathbf{Q}_{p}$ denote the ring of $p$-adic integers and the field of $p$-adic numbers, respectively.

For an associative ring $R$ with identity element and an $R$-module $M$ we denote by $R^{\times}$the group of all its invertible elements and by $M_{n}^{m}$ the $R$-module of all $m \times n$-matrices with entries in $M$; we put $M^{m}=M_{1}^{m}$ for simplicity. Sometimes an object with a superscript such as $G^{n}$ in Section 23 is used with a different meaning, but the distinction will be clear from the context. For $x \in R_{n}^{m}$ and an ideal $\mathfrak{a}$ of $R$ we write $x \prec \mathfrak{a}$ if all the entries of $x$ belong to $\mathfrak{a}$. (There is a variation of this; see §1.8.)

The transpose, determinant, and trace of a matrix $x$ are denoted by ${ }^{t} x, \operatorname{det}(x)$, and $\operatorname{tr}(x)$. The zero element of $R_{n}^{m}$ is denoted by $0_{n}^{m}$ or simply by 0 , and the identity element of $R_{n}^{n}$ by $1_{n}$ or simply by 1 . The size of a zero matrix block written simply 0 should be determined by the size of adjacent nonzero matrix blocks. We put $G L_{n}(R)=\left(R_{n}^{n}\right)^{\times}$, and

$$
S L_{n}(R)=\left\{\alpha \in G L_{n}(R) \mid \operatorname{det}(\alpha)=1\right\}
$$

if $R$ is commutative. If $x_{1}, \ldots, x_{r}$ are square matrices, $\operatorname{diag}\left[x_{1}, \ldots, x_{r}\right]$ denotes the matrix with $x_{1}, \ldots, x_{r}$ in the diagonal blocks and 0 in all other blocks. We shall be considering matrices $x$ with entries in a ring with an anti-automorphism $\rho$ (complex conjugation, for example), including the identity map. We then put $x^{*}={ }^{t} x^{\rho}$, and $\widehat{x}=\left(x^{*}\right)^{-1}$ if $x$ is square and invertible.

For a complex number or more generally for a complex matrix $\alpha$ we denote by $\operatorname{Re}(\alpha), \operatorname{Im}(\alpha)$, and $\bar{\alpha}$ the real part, the imaginary part, and the complex conjugate of $\alpha$. For complex hermitian matrices $x$ and $y$ we write $x>y$ and $y<x$ if $x-y$ is positive definite, and $x \geq y$ and $y \leq x$ if $x-y$ is nonnegative. For $r \in \mathbf{R}$ we denote by $[r]$ the largest integer $\leq r$.

Given a set $A$, the identity map of $A$ onto itself is denoted by $\mathrm{id}_{A}$ or $1_{A}$. To indicate that a union $X=\bigcup_{i \in I} Y_{i}$ is disjoint, we write $X=\bigsqcup_{i \in I} Y_{i}$. We understand that $\prod_{i=\alpha}^{\beta}=1$ and $\sum_{i=\alpha}^{\beta}=0$ if $\alpha>\beta$. For a finite set $X$ we denote by $\# X$ or $\#(X)$ the number of elements in $X$. If $H$ is a subgroup of a group $G$, we put $[G: H]=\#(G / H)$. However we use also the symbol $[K: F]$ for the degree
of an algebraic extension $K$ of a field $F$. The distinction will be clear from the context. By a Hecke character $\chi$ of a number field $K$ we mean a continuous Tvalued character of the idele group of $K$ trivial on $K^{\times}$, and denote by $\chi^{*}$ the ideal character associated with $\chi$. By a $C M$-field we mean a totally imaginary quadratic extension of a totally real algebraic number field.

## INTRODUCTION

Our ultimate aim is to prove several theorems of arithmeticity on the values of an Euler product $\mathcal{Z}(s)$ and an Eisenstein series $E(z, s)$ at certain critical points $s$. We take these $\mathcal{Z}$ and $E$ to be those of the types we treated in our previous book "Euler Products and Eisenstein Series," referred to as [S97] here. They are defined with respect to an algebraic group $G$, which is either symplectic or unitary. To illustrate the nature of our problems, let us take a CM-field $K$ and put

$$
\begin{equation*}
G(\varphi)=G^{\varphi}=\left\{\alpha \in G L_{n}(K) \mid \alpha \varphi \cdot{ }^{t} \alpha^{\rho}=\varphi\right\} \tag{0.1}
\end{equation*}
$$

where $\rho$ denotes complex conjugation and $\varphi$ is an element of $G L_{n}(K)$ such that ${ }^{t} \varphi^{\rho}=\varphi$. This group acts on a hermitian symmetric space which we write $\mathfrak{Z}^{\varphi}$. We shall often be interested in the special case where $\varphi$ takes the form

$$
\eta=\eta_{q}=\left[\begin{array}{cc}
0 & 1_{q}  \tag{0.2}\\
1_{q} & 0
\end{array}\right]
$$

In this case we write $\mathcal{H}_{q}$, or simply $\mathcal{H}$, instead of $\mathcal{Z}^{\varphi}$ for the symmetric space.
Given a congruence subgroup $\Gamma$ of $G$, a Hecke eigenform $\mathbf{f}$ of holomorphic type on $\mathbf{3}^{\varphi}$ with respect to $\Gamma$, and a Hecke character $\chi$ of $K$ of algebraic type, but not necessarily of finite order, we can construct a "twisted Euler product" $\mathcal{Z}(s, \mathbf{f}, \chi)$, whose generic Euler $p$-factor for each rational prime $p$ has degree $n[K: \mathbf{Q}]$. Then we shall eventually prove that

$$
\begin{equation*}
\mathcal{Z}\left(\sigma_{0}, \mathbf{f}, \chi\right) \in \pi^{\varepsilon} \mathfrak{q}\langle\mathbf{f}, \mathbf{f}\rangle \overline{\mathbf{Q}} \tag{0.3}
\end{equation*}
$$

for $\sigma_{0}$ in a certain finite subset of $2^{-1} \mathbf{Z}$ and $\overline{\mathbf{Q}}$-rational $\mathbf{f}$. Here $\langle\mathbf{f}, \mathbf{f}\rangle$ is the inner product defined in a canonical way; $\varepsilon$ is an integer determined by $\sigma_{0}$, the signature of $\varphi$, the weight of $\mathbf{f}$, and the archimedean factor of $\chi ; \mathfrak{q}$ is a certain "period symbol" determined by $\chi$ and $\varphi$. This is true for both isotropic and anisotropic $\varphi$, and even for a totally definite $\varphi$. In the simplest case in which $G=G^{\eta}$, we have $\mathfrak{q}=1$.

Clearly such a result requires the definition of $\overline{\mathbf{Q}}$-rationality of automorphic forms. If $G$ is of type $G^{\eta}$, then we can define the $\overline{\mathbf{Q}}$-rationality by the $\overline{\mathbf{Q}}$-rationality of the Fourier coefficients of a given automorphic form. If $[K: \mathbf{Q}]=2$, for example, then $\mathcal{H}$ is a tube domain of the form $\mathcal{H}=\left\{z \in \mathbf{C}_{q}^{q} \mid i\left(z^{*}-z\right)>0\right\}$, and a holomorphic automorphic form $f$ has an expansion

$$
\begin{equation*}
f(z)=\sum_{h} c(h) \exp (2 \pi i \cdot \operatorname{tr}(h z)) \quad(z \in \mathcal{H}) \tag{0.4}
\end{equation*}
$$

with $c(h) \in \mathbf{C}$, where $h$ runs over all nonnegative hermitian matrices belonging to a Z-lattice in $K_{q}^{q}$. Then for a subfield $M$ of $\mathbf{C}$ we say that $f$ is $M$-rational if $c(h) \in M$ for every $h$. This definition may look simplistic, but actually it is intrinsically the
right definition. To explain about this point, we first note that $\Gamma \backslash \mathfrak{Z}^{\varphi}$ has a structure of algebraic variety that has a natural model $W$ defined over $\overline{\mathbf{Q}}$. We call then a $\Gamma$ automorphic function (that is, $\Gamma$-invariant meromorphic function on $\mathcal{Z}^{\varphi}$ satisfying the cusp condition) $\overline{\mathbf{Q}}$-rational (or arithmetic) if it corresponds to a $\overline{\mathbf{Q}}$-rational function on $W$ in the sense of algebraic geometry. Now there are two basic facts:
(0.5) The value of a $\overline{\mathbf{Q}}$-rational automorphic function at any CM-point of $\mathbf{Z}^{\varphi}$, if finite, is algebraic.
(0.6) If $f$ and $g$ are $\overline{\mathbf{Q}}$-rational automorphic forms of the same weight, then $f / g$ is a $\overline{\mathbf{Q}}$-rational automorphic function.

Here a CM-point on $\mathcal{Z}^{\varphi}$ is defined to be the fixed point of a certain type of torus contained in $G$. If $G=G^{\eta}$ and $q=1$, then $\mathcal{H}$ is the standard upper half plane, and any point of $\mathcal{H}$ belonging to an imaginary quadratic field is a CM-point and vice versa. In such a special case, (0.5) and (0.6) follow from the classical theory of complex multiplication of elliptic modular functions. In more general cases, (0.5) was established by the author in the framework of canonical models. As for (0.6), it makes sense if $G=G^{\eta}$, and we can indeed give a proof, if nontrivial, of (0.6) in such a case. For $\varphi$ of a more general type, however, (0.6) is a meaningful statement only when we have defined the $\overline{\mathbf{Q}}$-rationality of automorphic forms. Thus it is one of our main tasks to define the notion so that (0.6) holds.

Turning our eyes to Eisenstein series, easily posable questions are as follows:
(i) Assuming that $E\left(z, \sigma_{0}\right)$ is finite, is $E\left(z, \sigma_{0}\right)$ as a function of $z$ holomorphic?
(ii) If that is so, is it $\overline{\mathbf{Q}}$-rational up to a well-defined constant?

Here we take meromorphic continuation of $E(z, s)$ to the whole $s$-plane, as we proved in [S97], into account. Every researcher of automorphic forms should be able to accept such questions naturally, since the answers to them for $G=S L_{2}(\mathbf{Q})$ are well-known and fundamental. There is a marked difference between the $\overline{\mathbf{Q}}$ rationality here and the arithmeticity of $\mathcal{Z}\left(\sigma_{0}\right)$, since the latter concerns $\sigma_{0}$ in an interval which can be large, while $E\left(z, \sigma_{0}\right)$ can be holomorphic in $z$ only at a single point $\sigma_{0}$. Now the interval, or rather the set of critical points belonging to the interval, is suggested by the functional equation for $\mathcal{Z}$, and we can find such a set even for $E(z, s)$ by means of its analytic properties. We cannot expect $E\left(z, \sigma_{0}\right)$ to be holomorphic in $z$ for every critical point $\sigma_{0}$ in the set. We should also note a classical example in the elliptic modular case:

$$
\begin{align*}
\left(-4 \pi^{2}\right)^{-1} \lim _{s \rightarrow+0} & \sum_{0 \neq(c, d) \in \mathbf{Z}^{2}}(c z+d)^{-2}|c z+d|^{-s}  \tag{0.7}\\
& =(4 \pi y)^{-1}-12^{-1}+2 \sum_{n=1}^{\infty}\left(\sum_{a \mid n} a\right) e^{2 \pi i n z} .
\end{align*}
$$

This is a nonholomorphic modular form of weight 2 , and there are similar nonholomorphic forms of weight $(n+3) / 2$ with respect to a congruence subgroup of $S p(n, \mathbf{Z})$. Therefore our next questions are:
(iii) What is the analytic nature of these $E\left(z, \sigma_{0}\right)$ ?
(iv) Can we still speak of the $\overline{\mathbf{Q}}$-rationality of such $E\left(z, \sigma_{0}\right)$ ?

One of the main purposes of this book is to answer these questions, which are not only meaningful by themselves, but also closely connected with the arithmeticity of
$\mathcal{Z}\left(\sigma_{0}\right)$. In fact, the answers to (iii) and (iv) are indispensable for the proof of (0.3) as we shall explain later, but first let us describe our answers.

We first define the notion of nearly holomorphic function on any complex manifold with a fixed Kähler structure. Without going into details in the general case, let us just say that a function on such a manifold $\mathcal{Z}$ is called nearly holomorphic if it is a polynomial of some functions $r_{1}, \ldots, r_{m}$ on $\mathfrak{Z}$, determined by the Kähler structure, over the ring of all holomorphic functions on $\mathfrak{Z}$. If $\mathfrak{Z}$ is the above $\mathcal{H}$ of tube type with a $G$-invariant Kähler structure, then the $r_{i}$ are the entries of $\left(z^{*}-z\right)^{-1}$, where $z$ is a variable matrix on $\mathcal{H}$. If $\mathcal{Z}$ is a hermitian symmetric space, there is also a characterization of such functions in terms of the Lie algebra of the transformation group on 3 .

Now we can naturally define nearly holomorphic automorphic forms by replacing holomorphy by near holomorphy in the definition of automorphic forms. If $G=G^{\eta}$, then such a form $f$ on $\mathcal{H}$ has an expansion

$$
\begin{equation*}
f(z)=\sum_{h} p_{h}\left(\left[\pi i\left(z^{*}-z\right)\right]^{-1}\right) \exp (2 \pi i \cdot \operatorname{tr}(h z)) \quad(z \in \mathcal{H}) \tag{0.8}
\end{equation*}
$$

where $\sum_{h}$ is the same as in (0.4) and $p_{h}(Y)$ is a polynomial function in the entries of $Y$ whose degree is less than a constant depending on $f$. We say that $f$ is $M$-rational if $p_{h}$ has all its coefficients in a field $M$ for every $h$. For example, the function of (0.7) is a Q-rational nearly holomorphic modular form. We shall show that $E\left(z, \sigma_{0}\right)$ is indeed nearly holomorphic and $\overline{\mathbf{Q}}$-rational in this sense, up to a constant, which is a power of $\pi$ if $G=G^{\eta}$. Moreover, here is a noteworthy consequence of our definition:
(0.9) If $f$ and $g$ are $\overline{\mathbf{Q}}$-rational nearly holomorphic automorphic forms of the same weight, then the value of $f / g$ at any $C M$-point of $\mathcal{H}$, if finite, is algebraic.

It should be noted that this is anything but a direct consequence of (0.6). Also, for a general type of $\varphi$ we cannot use (0.8). However, once we have the $\overline{\mathbf{Q}}$-rationality of holomorphic automorphic forms, we can at least define the $\overline{\mathbf{Q}}$-rationality of nearly holomorphic automorphic forms by property ( 0.9 ), though it is of course nontrivial to show that such a definition is indeed meaningful. So far we have taken $G$ to be unitary, but the symplectic case can be handled too; in fact it is similar to and easier than $G^{\eta}$, though the case of half-integral weight requires special consideration.

Having thus presented our problems in rough forms, we can now set our program as follows:
(1) We first define the $\overline{\mathbf{Q}}$-rationality of automorphic forms so that (0.6) holds.
(2) We define nearly holomorhic automorphic forms and their $\overline{\mathbf{Q}}$-rationality so that (0.9) holds.
(3) We prove the near holomorphy and $\overline{\mathbf{Q}}$-rationality of $E\left(z, \sigma_{0}\right)$ up to a power of $\pi$ in the easiest cases, namely, when $G$ is symplectic or of type $G^{\eta}$, and $E$ is defined with respect to a parabolic subgroup whose unipotent radical is a commutative group of translations on $\mathcal{H}$. Let us call such an $E$ a series of split type.
(4) We prove ( 0.3 ) by using the result of (3).
(5) Finally we prove the near holomorphy and $\overline{\mathbf{Q}}$-rationality of $E\left(z, \sigma_{0}\right)$ up to a well-defined constant in the most general case.

Let us now briefly describe the technical aspect of how these can be achieved. One important point is that certain differential operators on $\mathcal{H}$ are essential to (2) and (3). In the above we tacitly assumed that our automorphic forms are
scalar-valued, but in order to use differential operators effectively, it is necessary to consider vector-valued forms. If $[K: \mathbf{Q}]=2$ and $G=G\left(\eta_{q}\right)$, such a form is defined relative to a representation $\{\rho, X\}$ of a group

$$
\mathfrak{K}=\left\{(a, b) \in G L_{q}(\mathbf{C}) \times G L_{q}(\mathbf{C}) \mid \operatorname{det}(a)=\operatorname{det}(b)\right\},
$$

where $X$ is a finite-dimensional complex vector space and $\rho$ is a rational representation of $\mathfrak{K}$ into $G L(X)$. Put $T=\mathbf{C}_{q}^{q}$ and view it as a global holomorphic tangent space of $\mathcal{H}_{q}$; define a representation $\{\rho \otimes \tau, \operatorname{Hom}(T, X)\}$ of $\mathfrak{K}$ by $[(\rho \otimes \tau)(a, b) h](u)=\rho(a, b) h\left({ }^{t} a u b\right)$ for $(a, b) \in \mathfrak{K}, h \in \operatorname{Hom}(T, X)$, and $u \in T$. For a function $g: \mathcal{H} \rightarrow X$ we define $\operatorname{Hom}(T, X)$-valued function $D g$ and $D_{\rho} g$ on $\mathcal{H}$ by

$$
\begin{aligned}
& (D g)(u)=\sum_{i, j=1}^{q} u_{i j} \partial g / \partial z_{i j} \quad(u \in T) \\
& \left(D_{\rho} g\right)(z)=\rho(\Xi(z))^{-1} D[\rho(\Xi(z)) g(z)]
\end{aligned}
$$

where $z=\left(z_{i j}\right)_{i, j=1}^{q} \in \mathcal{H}$ and $\Xi: \mathcal{H} \rightarrow \mathfrak{K}$ is defined by $\Xi(z)=\left(i\left(\bar{z}-{ }^{t} z\right), i\left(z^{*}-z\right)\right)$. These can also be defined on $\mathfrak{Z}^{\varphi}$ for $\varphi$ of a general type. Then we can show that if $g$ is an automorphic form of weight $\rho$, then $D_{\rho} g$ is a form of weight $\rho \otimes \tau$. If $q=1$, then $\mathcal{H}$ is the standard upper half plane, $G^{\eta} \cap S L_{2}(K)=S L_{2}(\mathbf{Q}), \mathfrak{K}=\mathbf{C}^{\times}$, $\rho(a)=a^{k}$ for $a \in \mathbf{C}^{\times}$with $k \in \mathbf{Z}$, and $\Xi(z)=(2 y, 2 y)$ where $y=\operatorname{Im}(z)$; we can easily identify $D g$ with $\partial g / \partial z$, so that $D_{\rho} g=y^{-k}(\partial / \partial z)\left(y^{k} g\right)$, and $(\rho \otimes \tau)(a)=$ $a^{k+2}$. Thus $D_{\rho}$ is the well-known operator that sends a form of weight $k$ to a form of weight $k+2$.

Now iteration of operators of this type, such as $D_{\rho \otimes \tau} D_{\rho}$, produces an automorphic form with values in a representation space of $\mathfrak{K}$ of a large dimension if $q>1$, even if we start with $X=\mathbf{C}$. Decomposing the space into irreducible subspaces and looking particularly at the irreducible subspaces of dimension one, we can define a natural differential operator $\Delta$ that sends scalar-valued automorphic forms to scalar-valued forms of increased weight. The significance of these iterated operators and $\Delta$ are explained by the following fact, which is formulated only for $\Delta$ for simplicity:
(0.10) If $\Delta$ is of total degree $p$ in terms of $\partial / \partial z_{i j}$, then $\pi^{-p} \Delta$ preserves near holomorphy and $\overline{\mathbf{Q}}$-rationality.
If $G=G^{\eta}$, this can be derived from our definition in terms of expression (0.8). Now property ( 0.9 ), if true, would imply that for a $\overline{\mathbf{Q}}$-rational holomorphic automorphic forms $f$ and $g$ such that $\Delta f$ and $g$ have the same weight, the value of $\left(\pi^{-p} \Delta f\right) / g$ at any CM-point, if finite, is algebraic. This is highly nontrivial, and in fact we first prove this special case of (0.9), and derive the general case from that result.

As for problem (3), we first investigate the Fourier expansion of $E(z, s)$ of split type. In fact, this was done in [S97], but here we examine the behavior of the Fourier coefficients at a critical value of $s$. Employing their explicit forms, we find that $E(z, \sigma)$ is holomorphic in $z$ and $\overline{\mathbf{Q}}$-rational, or is of the type (0.7), if the weight of $E$ and $\sigma$ belong to certain special types. For a more general weight and a general $\sigma_{0}$, we prove that $c E\left(z, \sigma_{0}\right)=\Delta E^{\prime}(z, \sigma)$ with a suitable $\Delta$, a nonzero constant $c$, and a suitable $E^{\prime}$ belonging to those special types. Then ( 0.10 ) settles problem (3) for $E\left(z, \sigma_{0}\right)$.

To treat problems (4) and (5), let us now go back to the Euler product $\mathcal{Z}(s, \mathbf{f}, \chi)$ of $(0.3)$ on $G^{\varphi}$; we refer the reader to [S97] for its precise definition. We consider $G^{\psi}$ with $\psi=\operatorname{diag}[\varphi, \eta]$, where $\eta$ is as in (0.2). Then $G^{\varphi} \times G^{\eta}$ can be embedded
in $G^{\psi}$, and $G^{\psi}$ has a parabolic subgroup whose reductive factor is $G^{\varphi} \times G L_{q}(K)$. Given a suitable congruence subgroup $\Gamma^{\prime}$ of $G^{\psi}$, we can define an Eisenstein series $E(z, s ; \mathbf{f}, \chi)$ for $(z, s) \in \mathbf{3}^{\psi} \times \mathbf{C}$ with respect to that parabolic subgroup and the set of data ( $\mathbf{f}, \chi, \Gamma^{\prime}$ ). Now we easily see that $\operatorname{diag}[\psi,-\varphi]$ is equivalent to $\eta_{n+q}$, so that $G^{\psi} \times G^{\varphi}$ can be embedded into $G\left(\eta_{n+q}\right)$, and $\mathcal{Z}^{\psi} \times \mathfrak{Z}^{\varphi}$ can be embedded into $\mathcal{H}_{n+q}$. Pulling back an Eisenstein series on $\mathcal{H}_{n+q}$ of split type to $\boldsymbol{3}^{\psi} \times \boldsymbol{Z}^{\varphi}$, we obtain a function $H(z, w ; s)$ of $(z, w ; s) \in \mathfrak{Z}^{\psi} \times \mathfrak{Z}^{\varphi} \times \mathbf{C}$, with which we proved in [S97] an equality that takes the form

$$
\begin{equation*}
c(s) \mathcal{Z}(s, \mathbf{f}, \chi) E(z, s ; \mathbf{f}, \chi)=\Lambda(s) \int_{\Gamma \backslash 3^{\varphi}} H(z, w ; s) \mathbf{f}(w) \delta(w)^{m} \mathbf{d} w \tag{0.11}
\end{equation*}
$$

in the simplest case, where $\Gamma$ is a congruence subgroup of $G^{\varphi}, c$ is an easy product of gamma functions, $\Lambda$ is a product of some $L$-functions, $\mathbf{d} w$ is a $G^{\varphi}$-invariant measure on $\mathcal{Z}^{\varphi}$, and $\delta(w)^{m}$ is a factor, similar to $y^{k}$ in the one-dimensional case, that makes the integral meaningful. If $\psi=\varphi$, then (0.11) takes the form

$$
\begin{equation*}
c^{\prime}(s) \mathcal{Z}(s, \mathbf{f}, \chi) \mathbf{f}(z)=\Lambda^{\prime}(s) \int_{\Gamma \backslash \mathfrak{Z}^{\varphi}} H^{\prime}(z, w ; s) \mathbf{f}(w) \delta(w)^{m} \mathbf{d} w . \tag{0.12}
\end{equation*}
$$

We evaluate (0.11) and (0.12) at $s=\sigma_{0}$ for $\sigma_{0}$ belonging to a certain "critical set," and observe that $H\left(z, w ; \sigma_{0}\right)$ is nearly holomorphic in $(z, w) \in \mathcal{Z}^{\psi} \times \mathcal{Z}^{\varphi}$, and even $\overline{\mathbf{Q}}$-rational up to a power of $\pi$ and a factor $\mathfrak{q}$ as in (0.3). Then we can show that

$$
\Lambda\left(\sigma_{0}\right) H\left(z, w ; \sigma_{0}\right)=\pi^{\alpha} \mathfrak{q} \sum_{i} g_{i}(z) \overline{h_{i}(w)}
$$

with some $\alpha \in \mathbf{Z}$, and functions $g_{i}$ on $\mathcal{Z}^{\psi}$ and $h_{i}$ on $\mathfrak{Z}^{\varphi}$, which are nearly holomorphic and $\overline{\mathbf{Q}}$-rational. The same is true for $\Lambda^{\prime} H^{\prime}$; both $g_{i}$ and $h_{i}$ are defined on $\mathfrak{Z}^{\varphi}$ then. This fact applied to (0.12) produces a proportionality relation

$$
\mathcal{Z}\left(\sigma_{0}, \mathbf{f}, \chi\right) \in \pi^{\beta} \mathfrak{q}\left\langle\mathbf{p}^{\prime}, \mathbf{f}\right\rangle \overline{\mathbf{Q}}
$$

with some $\beta \in \mathbf{Z}$ and a $\overline{\mathbf{Q}}$-rational nearly holomorhic $\mathbf{p}^{\prime}$. Now we can show that $\mathcal{Z}(s, \mathbf{f}, \chi) \neq 0$ for $\operatorname{Re}(s)>3 q / 2$ if $G=G\left(\eta_{q}\right)$ and for $\operatorname{Re}(s)>n$ if $G=G^{\varphi}$ with $\varphi$ of a general type. There is one more crucial technical fact that we can replace $\mathbf{p}^{\prime}$ by a $\overline{\mathbf{Q}}$-rational holomorphic cusp form $\mathbf{p}$ that belongs to the same Hecke eigenvalues as $\mathbf{f}$. Choosing $\sigma_{0}$ so that $\mathcal{Z}\left(\sigma_{0}, \mathbf{f}, \chi\right) \neq 0$, we can show that $\langle\mathbf{p}, \mathbf{f}\rangle /\langle\mathbf{f}, \mathbf{f}\rangle \in \overline{\mathbf{Q}}$, and eventually ( 0.3 ) for $\sigma_{0}$ belonging to an appropriate set. Strictly speaking, ( 0.12 ) is true only under a consistency condition on ( $\mathbf{f}, \chi$ ), and the proof of ( 0.3 ) in the most general case is more complicated.

Next, we evaluate (0.11) at a critical $\sigma_{0}$ in a similar way, to find that

$$
\mathcal{Z}\left(s_{0}, \mathbf{f}, \chi\right) E\left(z, \sigma_{0} ; \mathbf{f}, \chi\right)=\pi^{\gamma} \mathfrak{q}\langle\mathbf{r}, \mathbf{f}\rangle g(z)
$$

with some $\overline{\mathbf{Q}}$-rational nearly holomorphic function $g$ on $\mathbf{Z}^{\psi}$ and some $\mathbf{r}$ of the same type as the above $\mathbf{p}$. Dividing this equality by $\langle\mathbf{f}, \mathbf{f}\rangle$ and employing (0.3), we obtain the desired near holomorphy and $\overline{\mathbf{Q}}$-rationality of $E\left(z, \sigma_{0} ; \mathbf{f}, \chi\right)$, which is the final main result of this book.

Since the title of each section can give a rough idea of its contents, we shall not describe them here for every section. However, there are some points which are not discussed in the above, and on which special attention may be paid. Let us note here some of the noteworthy aspects.
(A) As to the arithmeticity of automorphic forms, we stated only (0.6) as a basic requirement. However, there are other natural questions about arithmeticity whose answers become necessary in various applications. Let us mention here only a few facts we shall prove in this connection: (i) all automorphic forms are spanned by the $\overline{\mathbf{Q}}$-rational forms; (ii) the group action (defined relative to a fixed weight) preserves $\overline{\mathbf{Q}}$-rationality; (iii) in these statements $\overline{\mathbf{Q}}$ can be replaced by a smaller field such as $\mathbf{Q}$ or $\mathbf{Q}_{\mathrm{ab}}$ if the group and the weight are of special types.
(B) In Sections 19 through 25 we give a detailed treatment of $\mathcal{Z}(s, \mathbf{f}, \chi)$ and $E(z, s ; \mathbf{f}, \chi)$ in the symplectic case, as well as in the case $G=G^{\eta}$. These cases were mentioned but not discussed in detail in the previous book [S97]. Also, in the symplectic case we can define $\mathcal{Z}$ and $E$ even with respect to a half-integral weight, and we believe that the subject acquires the status of a complete theory only when that case is included. Therefore in this book we treat both integral and half-integral weights, and present the main results for both, though at a few points the details of the proof for a half-integral weight are referred to some papers of the author.
(C) We have spoken of a CM-point, which is naturally related to an abelian variety with complex multiplication. Thus it is necessary to view $\Gamma \backslash \mathfrak{Z}^{\varphi}$ as a space parametrizing a family of abelian varieties. This will be discussed in Sections 4 and 6. The topic was treated in [S98], but we prove here something which was not fully explained in that book. Namely, in Section 9, we prove the reciprocity-law for the value of an automorphic function at a CM-point, when $\Gamma \backslash \mathfrak{Z}^{\varphi}$ is associated with a PEL-type.
(D) In the elliptic modular case it is well-known that the space of all holomorphic modular forms is the direct sum of the space of cusp forms and the space of Eisenstein series. In Section 27 we prove several results of the same nature for symplectic and unitary groups. For example, we show that the orthogonal complement of the space of cusp forms in the space of all holomorphic automorphic forms is spanned by certain Eisenstein series, and the direct sum decomposition can be done $\overline{\mathbf{Q}}$-rationally. This will be proven for the weights with which the series are defined beyond the line of convergence.
(E) Though we are mainly interested in the higher-dimensional cases, in Section 18 we give an elementary theory of Eisenstein series in the Hilbert modular case, which leads to arithmeticity results on the critical values of an $L$-function of a CMfield. Also, in the Appendix we include some material of expository nature such as theta functions of a quadratic form and the estimate of the Fourier coefficients of a modular form. Many of them are well-known when the group is $S L_{2}(\mathbf{Q})$ or even $S p(n, \mathbf{Q})$ for some statements, but the researchers have often had difficulties in finding references for the results on a more advanced level. Therefore we have expended conscious efforts in treating such standard topics in a rather general setting.

## CHAPTER I

## AUTOMORPHIC FORMS AND FAMILIES OF ABELIAN VARIETIES

## 1. Algebraic preliminaries

1.1. The algebraic or Lie groups we treat in this book are symplectic and unitary, and the hermitian symmetric domains associated with them belong to Types A and C. Our methods are in fact applicable to groups and domains of other types, but it is naturally cumbersome to treat all cases. Therefore, in order to keep the book a reasonable length, we confine ourselves to those two types, though at some points we shall indicate that other cases can be handled in a similar way by citing relevant papers.

We take a basic field $F$ of characteristic different from 2 and a couple ( $K, \rho$ ) consisting of an $F$-algebra $K$ of rank $\leq 2$ and an $F$-linear automorphism $\rho$ of $K$ belonging to the following three types:
(I) $K=F$ and $\rho=\operatorname{id}_{F}$;
(II) $K$ is a quadratic extension of $F$ and $\rho$ is the generator of $\operatorname{Gal}(K / F)$;
(III) $K=F \times F$ and $(x, y)^{\rho}=(y, x)$ for $(x, y) \in F \times F$.

In our later discussion, objects of type (III) will appear as the localizations of the global objects of type (II).

Given left $K$-modules $V$ and $W$, we denote by $\operatorname{Hom}(W, V ; K)$ the set of all $K$ linear maps of $W$ into $V$. We then put $\operatorname{End}(V, K)=\operatorname{Hom}(V, V ; K), G L(V, K)=$ $\operatorname{End}(V, K)^{\times}$, and $S L(V, K)=\{\alpha \in G L(V, K) \mid \operatorname{det}(\alpha)=1\}$. We drop the letter $K$ if that is clear from the context. We let $\operatorname{Hom}(W, V)$ act on $W$ on the right; namely we denote by $w \alpha$ the image of $w \in W$ under $\alpha \in \operatorname{Hom}(W, V)$.

Let $V$ be a left $K$-module isomorphic to $K_{m}^{1}$. Given $\varepsilon= \pm 1$, by an $\varepsilon$-hermitian form on $V$ we understand an $F$-linear map $\varphi: V \times V \rightarrow K$ such that

$$
\begin{gather*}
\varphi(x, y)^{\rho}=\varepsilon \varphi(y, x)  \tag{1.1}\\
\varphi(a x, b y)=a \varphi(x, y) b^{\rho} \quad \text { for every } \quad a, b \in K . \tag{1.2}
\end{gather*}
$$

Assuming $\varphi$ to be nondegenerate, we define groups $G U(\varphi), U(\varphi)$, and $S U(\varphi)$ by

$$
\begin{align*}
G U(\varphi) & =\left\{\alpha \in G L(V, K) \mid \varphi(x \alpha, y \alpha)=\nu(\alpha) \varphi(x, y) \quad \text { with } \nu(\alpha) \in F^{\times}\right\}  \tag{1.3}\\
U(\varphi) & =\{\alpha \in G U(\varphi) \mid \nu(\alpha)=1\}, \quad S U(\varphi)=U(\varphi) \cap S L(V, K) \tag{1.4}
\end{align*}
$$

We call $\varphi$ isotropic if $\varphi(x, x)=0$ for some $x \in V, \neq 0$; we call $\varphi$ anisotropic if $\varphi(x, x)=0$ only for $x=0$.

Given $(V, \varphi)$ and another structure ( $V^{\prime}, \varphi^{\prime}$ ) of the same type, we denote by $(V, \varphi) \oplus\left(V^{\prime}, \varphi^{\prime}\right)$ the structure $(W, \psi)$ given by $W=V \oplus V^{\prime}$ and $\psi\left(x+x^{\prime}, y+y^{\prime}\right)=$
$\varphi(x, y)+\varphi^{\prime}\left(x^{\prime}, y^{\prime}\right)$ for $x, y \in V$ and $x^{\prime}, y^{\prime} \in V^{\prime}$. We then view $U(\varphi) \times U\left(\varphi^{\prime}\right)$ as a subgroup of $U(\psi)$ in an obvious way.
1.2. We shall often express various objects by matrices. To simplify our notation, for a matrix $x$ with entries in $K$ we put

$$
\begin{equation*}
x^{*}={ }^{t} x^{\rho}, \quad x^{-\rho}=\left(x^{\rho}\right)^{-1}, \quad \widehat{x}={ }^{t} x^{-\rho}, \tag{1.5}
\end{equation*}
$$

assuming $x$ to be square and invertible if necessary. Now let $V=K_{m}^{1}$ and $\varphi=$ $\varepsilon \varphi^{*} \in K_{m}^{m}$. Then we can define an $\varepsilon$-hermitian form $\varphi_{0}$ on $V$ by $\varphi_{0}(x, y)=x \varphi y^{*}$ for $x, y \in V$. In this setting we shall always write simply $\varphi$ for the form $\varphi_{0}$. Then we have
(1.7) $U(\varphi)=\left\{\alpha \in G L_{m}(K) \mid \alpha \varphi \alpha^{*}=\varphi\right\}, \quad S U(\varphi)=U(\varphi) \cap S L_{m}(K)$.

We shall often consider $U\left(\eta_{n}\right)$ with

$$
\eta_{n}=\left[\begin{array}{cc}
0 & -1_{n}  \tag{1.8}\\
1_{n} & 0
\end{array}\right] .
$$

Here we are taking $\varepsilon=-1$. In particular, if $K=F$, the group $U\left(\eta_{n}\right)$ is usually denoted by $S p(n, F)$. More generally, for a commutative ring $A$ with identity element we put

$$
\begin{gather*}
S p(n, A)=\left\{\left.\alpha \in G L_{2 n}(A)\right|^{t} \alpha \eta_{n} \alpha=\eta_{n}\right\}  \tag{1.9}\\
G p(n, A)=\left\{\left.\alpha \in G L_{2 n}(A)\right|^{t} \alpha \eta_{n} \alpha=\nu(\alpha) \eta_{n} \text { with } \nu(\alpha) \in A^{\times}\right\} . \tag{1.10}
\end{gather*}
$$

Notice that

$$
\begin{array}{cc}
\operatorname{det}(\alpha)=\nu(\alpha)^{n} & (\alpha \in G p(n, A)), \\
\operatorname{det}(\alpha) \operatorname{det}(\alpha)^{\rho}=\nu(\alpha)^{m} & (\alpha \in G U(\varphi)) . \tag{1.12}
\end{array}
$$

The latter formula is obvious. To prove (1.11), let $\alpha \in G p(n, A)$ and $\beta=$ $\operatorname{diag}\left[1_{n}, \nu(\alpha) 1_{n}\right]$. Then $\beta \in G p(n, A)$ and $\nu(\beta)=\nu(\alpha)$; thus $\beta^{-1} \alpha \in S p(n, A)$. It is well-known that $\operatorname{det}(S p(n, A))=1$, and hence $\operatorname{det}(\alpha)=\operatorname{det}(\beta)=\nu(\alpha)^{n}$, which is (1.11). In particular $G p(1, A)=G L_{2}(A)$ and $S p(1, A)=S L_{2}(A)$.

Let $\xi=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right] \in G L_{2 n}(K)$ with $a, b, c, d$ of size $n$ and let $s \in F^{\times}$. Then

$$
\begin{align*}
& \xi \in G U\left(\eta_{n}\right) \text { and } \nu(\xi)=s \Longleftrightarrow a^{*} d-c^{*} b=s 1_{n}, a^{*} c=c^{*} a, b^{*} d=d^{*} b,  \tag{1.13}\\
& \Longleftrightarrow a d^{*}-b c^{*}=s 1_{n}, a b^{*}=b a^{*}, c d^{*}=d c^{*} .
\end{align*}
$$

1.3. Lemma. (1) Let $A$ be a commutative ring with identity element. Let $x=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right] \in G L_{m+n}(A)$ and $x^{-1}=\left[\begin{array}{ll}e & f \\ g & h\end{array}\right]$ with $a, e$ of size $m$ and $d, h$ of size $n$. Then $\operatorname{det}(x) \operatorname{det}(e)=\operatorname{det}(d)$.
(2) If $\xi=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right] \in S U\left(\eta_{n}\right)$ with $a, b, c, d$ of size $n$, then $\operatorname{det}(a), \operatorname{det}(b)$, $\operatorname{det}(c)$, and $\operatorname{det}(d)$ all belong to $F$.
(3) Every element of $G U\left(\eta_{n}\right)$ is a diagonal matrix times an element of $S U\left(\eta_{n}\right)$.
(4) $G U(\varphi) /\left[F^{\times} U(\varphi)\right]$ is isomorphic to a subgroup of $F^{\times} /\left\{a^{2} \mid a \in F^{\times}\right\}$. Consequently if $\lambda$ is a homomorphism of $G U(\varphi)$ into a group whose kernel contains $F^{\times} U(\varphi)$, then $\lambda^{2}=1$.

Proof. For the proof of (1) and (2), see [S97, Lemmas 2.15 and 2.16]. To prove (3), let $\alpha \in G U\left(\eta_{n}\right), p=\nu(\alpha)^{-1}$, and $\beta=\operatorname{diag}\left[1_{n}, p 1_{n}\right] \alpha$. Then $\beta \in U\left(\eta_{n}\right)$.

This settles our problem if $K=F$, since $U\left(\eta_{n}\right)=S U\left(\eta_{n}\right)$; so assume $K \neq F$. Put $q=\operatorname{det}(\beta)$. Then $q q^{\rho}=1$ by (1.12), and hence $q=r / r^{\rho}$ with some $r \in K^{\times}$. Put $\gamma=\operatorname{diag}\left[a^{*}, a^{-1}\right]$ with any diagonal matrix $a$ such that $\operatorname{det}(a)=r$. Then $\gamma \in U\left(\eta_{n}\right)$ and $\operatorname{det}(\gamma \beta)=1$. This proves (3) when $K \neq F$. Finally, consider the homomorphism $\nu: G U(\varphi) \rightarrow F^{\times}$. We easily see that $F^{\times} U(\varphi)$ is the inverse image of $\left\{a^{2} \mid a \in F^{\times}\right\}$, and hence we obtain (4).
1.4. Let $F$ be the field of quotients of a Dedekind domain $\mathfrak{g}$. By a $\mathfrak{g}$-lattice in a finite-dimensional vector space $W$ over $F$ we mean a finitely generated $\mathfrak{g}$-module in $W$ that spans $W$ over $F$. Every fractional ideal in $F$ with respect to $\mathfrak{g}$ is a $\mathfrak{g}$-lattice in $F$ and vice versa; we call it a $\mathfrak{g}$-ideal. A $\mathfrak{g}$-ideal is called integral if it is contained in $\mathfrak{g}$.

We now assume that $F$ is an algebraic number field of finite degree. We denote by $\mathbf{a}$ and $\mathbf{h}$ the sets of archimedean primes and nonarchimedean primes of $F$; we put $\mathbf{v}=\mathbf{a} \cup \mathbf{h}$. Further we denote by $\mathfrak{g}$ the maximal order of $F$. For every $v \in \mathbf{v}$ we denote by $F_{v}$ the $v$-completion of $F$. In particular, for $v \in \mathbf{h}$ and a $\mathfrak{g}$-ideal $\mathfrak{a}$ we denote by $\mathfrak{a}_{v}$ the $v$-closure of $\mathfrak{a}$ in $F_{v}$, which coincides with the $\mathfrak{g}_{v}$-linear span of $\mathfrak{a}$ in $F_{v}$. We denote by $N(\mathfrak{a})$ and $N\left(\mathfrak{a}_{v}\right)$ the norm of $\mathfrak{a}$ and $\mathfrak{a}_{v}$ as usual. They are positive rational numbers with the standard multiplicative property such that $N(\mathfrak{a})=[\mathfrak{g}: \mathfrak{a}]$ if $\mathfrak{a}$ is integral and $N\left(\mathfrak{a}_{v}\right)=\left[\mathfrak{g}_{v}: \mathfrak{a}_{v}\right]$ if $\mathfrak{a}_{v}$ is integral.

Given a finite-dimensional vector space $X$ over $F$ and a $\mathfrak{g}$-lattice $L$ in $X$, we put $X_{v}=X \otimes_{F} F_{v}$ for every $v \in \mathbf{v}$, and denote by $L_{v}$ the $\mathfrak{g}_{v}$-linear span of $L$ in $X_{v}$ if $v \in \mathbf{h}$. Clearly $L_{v}$ is a $\mathfrak{g}_{v}$-lattice in $X_{v}$, and is the closure of $L$ in $X_{v}$. Notice also that every $\mathfrak{g}_{v}$-lattice in $X_{v}$ is an open compact subgroup of $X_{v}$.
1.5. Lemma. With $F$ and $X$ as above, let $L$ be an arbitrarily fixed $\mathfrak{g}$-lattice in $X$. Then the following assertions hold:
(1) If $M$ is a $\mathfrak{g}$-lattice in $X$, then $L_{v}=M_{v}$ for almost all $v$. Moreover, $L \subset M$ (resp. $L=M$ ) if $L_{v} \subset M_{v}\left(\right.$ resp. $L_{v}=M_{v}$ ) for every $v \in \mathbf{h}$.
(2) Given a $\mathfrak{g}_{v}$-lattice $N_{v}$ in $X_{v}$ for each $v \in \mathbf{h}$ such that $N_{v}=L_{v}$ for almost all $v$, there exists a $\mathfrak{g}$-lattice $M$ in $X$ such that $M_{v}=N_{v}$ for every $v \in \mathbf{h}$.

These assertions are well-known. For the proof, see [S97, Lemma 8.2].
1.6. Given an algebraic group $G$ over $F$, we denote by $G_{\mathbf{A}}$ the adelization of $G$ and by $G_{v}$ for $v \in \mathbf{v}$ the localization of $G$ at $v$. (The reader is referred to [S97, Section 8] for basic definitions and elementary facts on this topic.) We consider $G$ a subgroup of $G_{\mathbf{A}}$ as usual. In particular, $F_{\mathbf{A}}$ and $F_{\mathbf{A}}^{\times}$denote the adele ring and the idele group of $F$, respectively. The archimedean and nonarchimedean factors of $G_{\mathbf{A}}$ are denoted by $G_{\mathbf{a}}$ and $G_{\mathbf{h}}$. Namely $G_{\mathbf{a}}=\prod_{v \in \mathbf{a}} G_{v}$ and $G_{\mathbf{h}}=G_{\mathbf{A}} \cap \prod_{v \in \mathbf{h}} G_{v}$. For $x \in G_{\mathbf{A}}$ we denote by $x_{\mathbf{a}}, x_{\mathbf{h}}$, and $x_{v}$ the projection of $x$ to $G_{\mathbf{a}}, G_{\mathbf{h}}$, and $G_{v}$, respectively. If $G \subset G L(V)$ with a vector space $V$ over $F$, then for $\alpha \in G_{\mathbf{A}}$ and a $\mathfrak{g}$-lattice $L$ in $V$, we denote by $L \alpha$ the $\mathfrak{g}$-lattice in $V$ determined by $(L \alpha)_{v}=L_{v} \alpha_{v}$ for every $v \in \mathbf{h}$. The existence of such a lattice $L \alpha$ is guaranteed by Lemma 1.5 (2). In particular, for $x \in F_{\mathbf{A}}^{\times}$we denote by $x \mathfrak{g}$ the fractional ideal such that $(x \mathfrak{g})_{v}=x_{v} \mathfrak{g}_{v}$. Also we put $|x|_{\mathbf{A}}=\prod_{v \in \mathbf{v}}\left|x_{v}\right|_{v}$, where $\left|\left.\right|_{v}\right.$ is the normalized valuation at $v$. To emphasize that this is defined on $F_{\mathbf{A}}^{\times}$, we shall also write $|x|_{F}$ for $|x|_{\mathbf{A}}$.

Given algebraic groups $G$ and $G^{\prime}$ over $F$ and an $F$-rational homomorphism $f$ of $G$ into $G^{\prime}$, we can extend $f$ naturally to a homomorphism of $G_{\mathbf{A}}$ to $G_{\mathbf{A}}^{\prime}$, which we shall denote by the same letter $f$. For example, we employ $\operatorname{Tr}_{F^{\prime} / F}$ even for the
map of $F_{\mathbf{A}}^{\prime}$ into $F_{\mathbf{A}}$ derived from the map $\operatorname{Tr}_{F^{\prime} / F}: F^{\prime} \rightarrow F$ when $F^{\prime}$ is an algebraic extension of $F$.

Let $W$ be a finite-dimensional vector space over $F$, and $L$ a $\mathfrak{g}$-lattice in $W$. Taking an element $a$ of $W_{\mathbf{h}}$, we define a function $\lambda$ on $W_{\mathbf{h}}$ by $\lambda(x)=\prod_{v \in \mathbf{h}} \lambda_{v}\left(x_{v}\right)$ for $x \in W_{\mathbf{h}}$, where $\lambda_{v}$ is the characteristic function of the coset $L_{v}+a_{v}$. We then denote by $\mathcal{S}\left(W_{\mathbf{h}}\right)$ the vector space of all finite $\mathbf{C}$-linear combinations of such functions $\lambda$ for all possible choices of ( $L, a$ ). This is called the Schwartz-Bruhat space of $W_{\mathbf{h}}$. We view every $\ell \in \mathcal{S}\left(W_{\mathbf{h}}\right)$ as a function on $W_{\mathbf{A}}$ by putting $\ell(x)=\ell\left(x_{\mathbf{h}}\right)$ for $x \in W_{\mathbf{A}}$. In particular, $\ell(\xi)$ is meaningful for every $\xi \in W$. We can easily see that the restriction of the elements of $\mathcal{S}\left(W_{\mathbf{h}}\right)$ to $W$ gives an isomorphism of $\mathcal{S}\left(W_{\mathbf{h}}\right)$ onto the set of all finite $\mathbf{C}$-linear combinations of functions, each of which is the characteristic function of a coset of $W$ modulo a $\mathfrak{g}$-lattice in $W$. This is because $W_{\mathbf{A}}=W+Y$ with $Y=\left\{y \in W_{\mathbf{A}} \mid y_{\mathbf{h}} \in \prod_{v \in \mathbf{h}} L_{v}\right\}$ for any fixed $\mathfrak{g}$-lattice $L$.

We now put

$$
\begin{equation*}
\mathbf{e}(z)=e^{2 \pi i z} \quad(z \in \mathbf{C}) \tag{1.14}
\end{equation*}
$$

and define characters $\mathbf{e}_{\mathbf{A}}: F_{\mathbf{A}} \rightarrow \mathbf{T}$ and $\mathbf{e}_{v}: F_{v} \rightarrow \mathbf{T}$ for each $v \in \mathbf{v}$ as follows: if $v \in \mathbf{a}$, then $\mathbf{e}_{v}(x)=\mathbf{e}(x)$ for real $v$ and $\mathbf{e}_{v}(x)=\mathbf{e}(x+\bar{x})$ for imaginary $v$; if $v \in \mathbf{h}$ and $p$ is the rational prime divisible by $v$, then $\mathbf{e}_{v}(x)=\mathbf{e}_{p}\left(\operatorname{Tr}_{F_{v} / \mathbf{Q}_{p}}(x)\right)$, where $\mathbf{e}_{p}(z)=\mathbf{e}(-y)$ with $y \in \bigcup_{m=1}^{\infty} p^{-m} \mathbf{Z}$ such that $z-y \in \mathbf{Z}_{p}$. We then put $\mathbf{e}_{\mathbf{A}}(x)=\prod_{v \in \mathbf{v}} \mathbf{e}_{v}\left(x_{v}\right), \mathbf{e}_{\mathbf{h}}(x)=\mathbf{e}_{\mathbf{A}}\left(x_{\mathbf{h}}\right)$, and $\mathbf{e}_{\mathbf{a}}(x)=\mathbf{e}_{\mathbf{A}}\left(x_{\mathbf{a}}\right)$ for $x \in F_{\mathbf{A}}$. We note here a basic property of $\mathbf{e}_{v}$ :

$$
\mathfrak{d}(F / \mathbf{Q})_{v}^{-1}=\left\{x \in F_{v} \mid \mathbf{e}_{v}(x y)=1 \text { for every } y \in \mathfrak{g}_{v}\right\} \quad(v \in \mathbf{h})
$$

where $\mathfrak{d}(F / \mathbf{Q})$ denotes the different of $F$ relative to $\mathbf{Q}$.
We insert here an easy fact as an application of Lemma 1.3 (4):
$\left.{ }^{*}\right)$ For every $\alpha \in G U(\varphi)_{\mathbf{A}}$ the map $x \mapsto \alpha x \alpha^{-1}$ of $U(\varphi)_{\mathbf{A}}$ onto itself leaves any fixed Haar measure of $U(\varphi)_{\mathbf{A}}$ invariant.

Indeed, let $\mu$ be a Haar measure of $U(\varphi)_{v}$ for a fixed $v \in \mathbf{v}$. Then, for $\alpha \in$ $G U(\varphi)_{v}$ we have $\mu\left(\alpha X \alpha^{-1}\right)=\lambda(\alpha) \mu(X)$ for every measurable set $X$ in $U(\varphi)_{v}$ with a positive real number $\lambda(\alpha)$. Clearly $\lambda$ is a homomorphism of $G U(\varphi)_{v}$ into $\mathbf{R}^{\times}$ and $\lambda\left(F_{v}^{\times}\right)=1$. Also, $\lambda\left(U(\varphi)_{v}\right)=1$ by [S97, Proposition 8.13 (1)]. Therefore, by Lemma 1.3 (4) we have $\lambda(\alpha)=1$, since $\lambda(\alpha)>0$.
1.7. Let $A$ be a principal ideal domain and $F$ the field of quotients of $A$. We call an element $X$ of $A_{n}^{m}$ primitive if $\operatorname{rank}(X)=\operatorname{Min}(m, n)$ and the elementary divisors of $X$ are all equal to $A$. If $m=n$, clearly $X$ is primitive if and only if $X \in G L_{n}(A)$. If $m<n$ (resp. $m>n$ ), then $X$ is primitive if and only if $X$ is the first $m$ rows (resp. $n$ columns) of an element of $G L_{n}(A)$ (resp. $G L_{m}(A)$ ). (For these and other properties of primitive matrices, see [S97, Lemmas 3.3 and 3.4].)

Given $x \in F_{n}^{m}$, we can find $c \in A_{n}^{m}$ and $d \in G L_{m}(F) \cap A_{m}^{m}$ such that $\left[\begin{array}{cc}c & d\end{array}\right]$ is primitive and $x=d^{-1} c$. We then call the last equality a (left) reduced expression for $x$, and define an integral ideal $\nu_{0}(x)$ by

$$
\begin{equation*}
\nu_{0}(x)=\operatorname{det}(d) A \tag{1.15}
\end{equation*}
$$

This is independent of the choice of $c$ and $d$. We call $\nu_{0}(x)$ the denominator ideal of $x$. We easily see that $\nu_{0}(x+a)=\nu_{0}(x)$ if $a \prec A$ (see Notation). For these and other properties of the symbol $\nu_{0}$ see [S97, Proposition 3.6, §3.7, and Lemma 3.8].

We now consider our setting to be that of $\S 1.4$, with an algebraic number field as $F$. Given $x \in\left(F_{v}\right)_{n}^{m}$ with $v \in \mathbf{h}$, we can naturally define $\nu_{0}(x)$ to be an integral $\mathfrak{g}_{v}$-ideal, taking $\mathfrak{g}_{v}$ to be $A$. Then we put

$$
\begin{equation*}
\nu(x)=N\left(\nu_{0}(x)\right)=\left[\mathfrak{g}_{v}: \nu_{0}(x)\right] . \tag{1.16}
\end{equation*}
$$

If $x=d^{-1} c$ is a left reduced expression for $x$, then $\nu(x)=|\operatorname{det}(d)|_{v}^{-1}$. Moreover, if $u x v=\left[\begin{array}{ll}a & 0 \\ 0 & 0\end{array}\right]$ with $u \in G L_{m}\left(\mathfrak{g}_{v}\right), v \in G L_{n}\left(\mathfrak{g}_{v}\right)$, and $a=\operatorname{diag}\left[a_{1}, \ldots, a_{r}\right]$, then $\nu(x)^{-1}$ is the product of $\left|a_{i}\right|_{v}$ for all $i$ such that $a_{i} \notin \mathfrak{g}_{v}$ (see [S97, Lemma 3.8 (2)]).
1.8. Take our setting to be the same as in Cases I and II in §1.1 with an algebraic number field as $F$. Namely, $K=F$ or $K$ is a quadratic extension of $F$. We denote by $\mathfrak{r}$ the ring of algebraic integers in $K$ and by $\mathbf{k}$ the set of all nonarchimedean primes of $K$. (Thus $\mathfrak{r}=\mathfrak{g}$ and $\mathbf{k}=\mathbf{h}$ if $K=F$.) Given $v \in \mathbf{k}$, an $\mathfrak{r}_{v}$-ideal $\mathfrak{a}$, and a matrix $x$ with entries in $K_{v}$, we write $x \prec \mathfrak{a}$ if all the entries of $x$ belong to $\mathfrak{a}$. Similarly, for a matrix $y$ with entries in $K_{\mathbf{A}}$ and an $\mathfrak{r}$-ideal $\mathfrak{b}$, we write $y \prec \mathfrak{b}$ if all the entries of $y_{v}$ belong to $\mathfrak{b}_{v}$ for every $v \in \mathbf{k}$.

Take two positive integers $m$ and $n$. For $x=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right] \in G L_{m+n}(K)_{\mathbf{A}}$ with $a \in\left(K_{\mathbf{A}}\right)_{m}^{m}$ and $d \in\left(K_{\mathbf{A}}\right)_{n}^{n}$, write $a=a_{x}, b=b_{x}, c=c_{x}$, and $d=d_{x}$. With fixed $\mathfrak{r}$-ideals $\mathfrak{y}$ and $\mathfrak{z}$ such that $\mathfrak{y z} \subset \mathfrak{r}$, we put

$$
\begin{align*}
C[\mathfrak{y}, \mathfrak{z}]=\left\{x \in G L_{m+n}(K)_{\mathbf{A}} \mid\right. & \operatorname{det}(x)_{\mathbf{h}} \in \prod_{v \in \mathbf{h}} \mathfrak{r}_{v}^{\times},  \tag{1.17}\\
& \left.a_{x} \prec \mathfrak{r}, b_{x} \prec \mathfrak{y}, c_{x} \prec \mathfrak{z}, d_{x} \prec \mathfrak{r}\right\} .
\end{align*}
$$

We easily see that this is a subgroup of $G L_{m+n}(K)_{\mathbf{A}}$ (see [S97, §9.1]). We also note that if $x \in C[\mathfrak{y}, \mathfrak{z}]$ and $v \mid \mathfrak{y z}$, then $\left(\operatorname{det}\left(a_{x}\right) \operatorname{det}\left(d_{x}\right)-\operatorname{det}(x)\right)_{v} \in \mathfrak{y}_{v} \mathfrak{z} v$, and hence we see that
(1.18) The map $x \mapsto\left(\left(a_{x}, b_{x}\right)_{v}\right)_{v \mid \mathfrak{n} \mathfrak{z}}$ defines a homomorphism of $C[\mathfrak{y}, \mathfrak{z}]$ into

$$
\prod_{v \mid \mathfrak{\mathfrak { n } _ { \mathfrak { z } }}}\left[G L_{m}\left(\mathfrak{r}_{v} / \mathfrak{y}_{v \mathfrak{z}_{v}}\right) \times G L_{n}\left(\mathfrak{r}_{v} / \mathfrak{y}_{v \mathfrak{z}_{v}}\right)\right] .
$$

1.9. Lemma. Define a subgroup $P^{(m, n)}$ of $G L_{m+n}(K)$ by

$$
P^{(m, n)}=\left\{x \in G L_{m+n}(K) \mid c_{x}=0\right\}
$$

Let $C$ denote the group $C[\mathfrak{y}, \mathfrak{z}]$ of (1.16). Then

$$
\begin{aligned}
P_{\mathbf{A}}^{(m, n)} C=\left\{x \in G L_{m+n}(K)_{\mathbf{A}} \mid\right. & \left(d_{x}\right)_{v} \in G L_{n}\left(K_{v}\right) \text { and } \\
& \left.\left(d_{x}^{-1} c_{x}\right)_{v} \prec \mathfrak{z}_{v} \text { for every } v \mid \mathfrak{y z}\right\} .
\end{aligned}
$$

Moreover, assuming $m=n$, let $G$ denote $G U\left(\eta_{n}\right), U\left(\eta_{n}\right)$, or $S U\left(\eta_{n}\right)$ with $\eta_{n}$ of (1.8); put $P=G \cap P^{(n, n)}$, and $D=G_{\mathbf{A}} \cap C$. Then

$$
\begin{aligned}
P_{\mathbf{A}} D & =G_{\mathbf{A}} \cap P_{\mathbf{A}}^{(n, n)} C \\
& =\left\{x \in G_{\mathbf{A}} \mid\left(d_{x}\right)_{v} \in G L_{n}(K)_{v} \text { and }\left(d_{x}^{-1} c_{x}\right)_{v} \prec \mathfrak{z}_{v} \text { for every } v \mid \mathfrak{n z}\right\} .
\end{aligned}
$$

In particular, $G L_{m+n}(K)_{\mathbf{A}}=P_{\mathbf{A}}^{(m, n)} C$ and $G_{\mathbf{A}}=P_{\mathbf{A}} D$ if $\mathfrak{y z}=\mathfrak{r}$.
Proof. The assertions for $G L_{m+n}(K)$ and $U\left(\eta_{n}\right)$ are proved in [S97, Lemma 9.2]. Combining the result for $U\left(\eta_{n}\right)$ with [S97, Lemma 9.10 (2)], we obtain the assertion for $S U\left(\eta_{n}\right)$. As for $G U\left(\eta_{n}\right)$, let $x \in G U\left(\eta_{n}\right)_{\mathbf{A}}, p=\operatorname{diag}\left[1_{n}, \nu(\alpha) 1_{n}\right]$, and $y=p^{-1} x$. Then $y \in U\left(\eta_{n}\right)$. Suppose $\left(d_{x}\right)_{v} \in G L_{n}(K)_{v}$ and $\left(d_{x}^{-1} c_{x}\right)_{v} \prec \mathfrak{z}_{v}$ for every $v \mid \mathfrak{n z}$. Then we easily see that $y$ satisfies the same conditions, and so
$y \in\left(U\left(\eta_{n}\right)_{\mathbf{A}} \cap P_{\mathbf{A}}^{(n, n)}\right)\left(U\left(\eta_{n}\right)_{\mathbf{A}} \cap C\right)$. Since $p$ is diagonal, we obtain the desired result for $G U\left(\eta_{n}\right)$.
1.10. The notation being as in $\S 1.8$, put $C[\mathfrak{z}]=C\left[\mathfrak{z}^{-1}, \mathfrak{z}\right]$. Then $G L_{m+n}(K)_{\mathbf{A}}=$ $P_{\mathbf{A}}^{(m, n)} C[\mathfrak{z}]$ by the above lemma. Therefore every element $x$ of $G L_{m+n}(K)_{\mathbf{A}}$ can be written $x=y z$ with $y \in P_{\mathbf{A}}^{(m, n)}$ and $z \in C[\mathfrak{z}]$. We then define an $\mathfrak{r}$-ideal il $\mathrm{il}_{\mathfrak{z}}(x)$ by

$$
\begin{equation*}
\mathrm{il}_{\mathfrak{z}}(x)=\operatorname{det}\left(d_{y}\right) \mathfrak{r} \tag{1.19}
\end{equation*}
$$

where the right-hand side is an $\mathfrak{r}$-ideal defined in $\S 1.6$. We easily see that this is well-defined. (But it depends on $m, n$, and $\mathfrak{z}$.)

Next, given $x \in\left(K_{\mathbf{A}}\right)_{m}^{n}$, we define an integral $\mathfrak{r}$-ideal $\nu_{0}(x)$ and a positive integer $\nu(x)$ by

$$
\begin{gather*}
\nu_{0}(x)_{v}=\nu_{0}\left(x_{v}\right) \text { for every } \quad v \in \mathbf{k}  \tag{1.20}\\
\nu(x)=N\left(\nu_{0}(x)\right) \tag{1.21}
\end{gather*}
$$

Here $\nu_{0}\left(x_{v}\right)$ is defined by (1.15). In other words, take $g \in G L_{n}(K)_{\mathbf{A}}$ and $h \in$ $\left(K_{\mathbf{A}}\right)_{m}^{n}$ so that $x_{v}=g_{v}^{-1} h_{v}$ is a reduced expression for every $v \in \mathbf{k}$. Then $\nu_{0}(x)=$ $\operatorname{det}(g) \mathbf{r}$.
1.11. Lemma. For $x \in G L_{m+n}(K)_{\mathbf{A}}$ the following assertions hold.
(1) If $\alpha=\operatorname{diag}\left[1_{m}, \kappa 1_{n}\right]$ with $\kappa \in K_{\mathbf{A}}^{\times}$, then $\mathrm{il}_{\kappa \mathfrak{z}}\left(\alpha x \alpha^{-1}\right)=\mathrm{il}_{\mathfrak{z}}(x)$.
(2) If $d_{x} \in G L_{n}(K)_{\mathbf{A}}$ and $\mathfrak{z}=\mu \mathbf{r}$ with $\mu \in K_{\mathbf{A}}^{\times}$, then $\operatorname{det}\left(d_{x}\right) \mathrm{il}_{\mathfrak{z}}(x)^{-1}=$ $\nu_{0}\left(\mu^{-1} d_{x}^{-1} c_{x}\right)$.
(3) If $x \in G L_{m+n}(K) \cap P_{\mathbf{A}}^{(m . n)} D[\mathfrak{n}, \mathfrak{z}]$ and $\mathfrak{y z} \neq \mathfrak{r}$, then $\operatorname{det}\left(d_{x}\right) \neq 0$ and $\operatorname{det}\left(d_{x}\right)$ . $\mathrm{il}_{\mathfrak{z}}(x)^{-1}$ is prime to $\mathfrak{y z}$.

For the proof see [S97, Lemma 9.4].
1.12. By a $C M$-field we mean a totally imaginary quadratic extension of a totally real algebraic number field of finite degree. Given a CM-field $K$ and an absolute equivalence class $\tau$ of representations of $K$ by complex matrices, we call ( $K, \tau$ ) a CM-type if the direct sum of $\tau$ and its complex conjugate is equivalent to the regular representation of $K$ over $\mathbf{Q}$. If $(K, \tau)$ is a CM-type and $[K: \mathbf{Q}]=2 n$, then $\tau$ is the class of $\operatorname{diag}\left[\tau_{1}, \ldots, \tau_{n}\right]$ with $n$ isomorphic embeddings $\tau_{i}$ of $K$ into $\mathbf{C}$ such that $\left\{\tau_{1}, \ldots, \tau_{n}, \tau_{1} \omega, \ldots, \tau_{n} \omega\right\}$ is exactly the set of all embeddings of $K$ into C, where $\omega$ denotes complex conjugation. In this setting we write $\tau=\left\{\tau_{i}\right\}_{i=1}^{n}$. If $F$ is the totally real field over which $K$ is quadratic and $\rho$ is the generator of $\operatorname{Gal}(K / F)$, then we have $\tau_{i} \omega=\rho \tau_{i}$, because of the following easy fact:
(1.22) If $\sigma$ is an isomorphism of $K$ onto a subfield of $\mathbf{C}$, then $x^{\rho \sigma}$ is the complex conjugate of $x^{\sigma}$ for every $x \in K$.

Therefore, if $X$ is a matrix with entries in $K$, then $\left(X^{*}\right)^{\sigma}={ }^{t} \overline{\left(X^{\sigma}\right)}$, where the bar is complex conjugation. Thus putting $Y^{*}={ }^{t} \bar{Y}$ for a complex matrix $Y$, we have $\left(X^{*}\right)^{\sigma}=\left(X^{\sigma}\right)^{*}$.

Given $(K, \tau)$ as above, let $K^{\prime}$ be the field generated over $\mathbf{Q}$ by $\sum_{i=1}^{n} a^{\tau_{i}}$ for all $a \in K$. We call $K^{\prime}$ the reflex field of $(K, \tau)$. It can easily be shown that $K^{\prime}$ is a CM-field and contains $\prod_{i=1}^{n} a^{\tau_{i}}$ for all $a \in K$ (see [S98, pp.62-63, p.122, Lemma 18.2]).

## 2. Polarized abelian varieties

In this section we review some basic facts on polarized abelian varieties defined over a subfield of $\mathbf{C}$. The reader can find more detailed treatments, as well as further references, in [W58] and [S98].
2.1. By an algebraic variety we understand an affine or a projective variety which is absolutely irreducible, defined in an affine or a projective space with a fixed coordinate system. If $V$ is an algebraic variety of dimension $n$, by a divisor of $V$ we mean a finite Z-linear combination of subvarieties of $V$ of dimension $n-1$.

Suppose now $V$ is an algebraic variety defined over a subfield of $\mathbf{C}$. Then by the same symbol $V$ we mean the point set consisting of all the points with coordinates in $\mathbf{C}$ satisfying the defining equations for $V$. If $V$ is nonsingular, then $V$ has a natural structure of a complex manifold.

By an abelian variety we uderstand a projective algebraic variety $A$ with a group structure such that the map $(x, y) \mapsto x+y$ of $A \times A$ into $A$ and also the map $x \mapsto-x$ of $A$ into $A$ are both rational maps defined everywhere. We use the additive notation, since such a group is always commutative. Such an $A$ must be nonsingular. We say that an abelian variety $A$ is defined over a field $k$ if the variety $A$ and these maps are defined over $k$. Thus, we speak of an abelian variety $A$ defined over a field $k$ always in this sense.

Let $A$ and $B$ be two abelian varieties. By a homomorphism of $A$ into $B$, or an endomorphism when $A=B$, we understand a rational map of $A$ into $B$ that is a group homomorphism. If such a map is birational, then it must be biregular, and we call it an isomorphism, or an automorphism when $A=B$. We denote by $\operatorname{Hom}(A, B)$ the set of all homomorphisms of $A$ into $B$, defined over any extension of a given field of definition for $A$ and $B$, and put $\operatorname{End}(A)=\operatorname{Hom}(A, A)$. We put also $\operatorname{Hom}_{\mathbf{Q}}(A, B)=\operatorname{Hom}(A, B) \otimes_{\mathbf{z}} \mathbf{Q}$ and $\operatorname{End}_{\mathbf{Q}}(A)=\operatorname{End}(A) \otimes_{\mathbf{z}} \mathbf{Q}$. An element of $\operatorname{Hom}(A, B)$ is called an isogeny if $A$ and $B$ have the same dimension and $\operatorname{Ker}(\lambda)$ is finite.
2.2. By a lattice in a finite-dimensional vector space $W$ over $\mathbf{R}$ we understand a discrete subgroup of $W$ that spans $W$ over $\mathbf{R}$. If $W$ is of dimension $m$, then a discrete subgroup $D$ of $W$ is a lattice in $W$ if and only if $D$ is isomorphic to $\mathbf{Z}^{m}$. (The distinction of a lattice in this sense from a $\mathfrak{g}$-lattice in the sense of $\S 1.4$ will be clear from the context, since the latter is never defined in a real vector space.)

Let $D$ be a lattice in $\mathbf{C}^{n}$. Then $D$ is isomorphic to $\mathbf{Z}^{2 n}$. An $\mathbf{R}$-valued $\mathbf{R}$-bilinear form $E: \mathbf{C}^{n} \times \mathbf{C}^{n} \rightarrow \mathbf{R}$ is called a Riemann form on $\mathbf{C}^{n}$ relative to $D$, or simply a Riemann form on $\mathbf{C}^{n} / D$, if it satisfies the following conditions:
(2.1) $E(D, D) \subset \mathbf{Z}$.
(2.2) $E(x, y)=-E(y, x)$.
(2.3) The form $(x, y) \mapsto E(x, i y)$ is symmeric and positive definite.

Usually a Riemann form is defined with nonnegativity instead of positive definiteness in the last condition. In the present book, however, we always consider positive definite Riemann forms, and so we take the above (2.1-3) to be the conditions for a Riemann form.

Now, given an abelian variety defined over a subfield of $\mathbf{C}$, there is always a complex-analytic biregular map, which is also a group isomorphism, of $A$ onto a complex torus. Conversely, it is a well-known fact that a complex torus $\mathbf{C}^{n} / D$ is
isomorphic to an abelian variety in that sense if and only if it has a Riemann form. If $\xi$ is such an isomorphism of $\mathbf{C}^{n} / D$ onto an abelian variety $A$, we can view $\xi$ as a homomorphism of $\mathbf{C}^{n}$ onto $A$ with kernel $D$. We then call $\left(\mathbf{C}^{n} / D, \xi\right)$ an analytic coordinate system of $A$.
2.3. Let $A$ and $\left(\mathbf{C}^{n} / D, \xi\right)$ be as above. Given a Riemann form $E$ on $\mathbf{C}^{n} / D$, put $H(u, v)=E(u, i v)+i E(u, v)$ for $u, v \in \mathbf{C}^{n}$. Then we can show that there is a nonzero holomorphic function $f$ on $\mathbf{C}^{n}$ such that

$$
\begin{equation*}
f(u+\ell)=f(u) \psi(\ell) \exp (\pi \cdot H(\ell, u+(\ell / 2))) \quad \text { for every } \quad \ell \in D \tag{2.4}
\end{equation*}
$$

with a map $\psi: D \rightarrow \mathbf{T}$. Then the zeros of $f$ define a divisor on $\mathbf{C}^{n} / D$, which we write $\xi^{-1}(X)$ with a divisor $X$ on $A$. In this situation we say that $E$ corresponds to $X$ or that $X$ determines $E$ with respect to $\xi$, since it can be shown that $E$ is unique for $X$, though $X$ is not uniquely determined by $E$.

Now we call a divisor on $A$ ample if it determines a Riemann form on $\mathbf{C}^{n} / D$. By a polarization of $A$ we mean a nonempty maximal set $\mathcal{C}$ of ample divisors of $A$ with the following property: if $X, Y \in \mathcal{C}$, then there exist positive integers $\ell, m$ such that $\ell X$ and $m Y$ determine the same Riemann form. By a polarized abelian variety we mean a structure $(A, \mathcal{C})$ formed by an abelian variety $A$ and a polarization $\mathcal{C}$ of $A$. We shall often denote $(A, \mathcal{C})$ by a single letter $\mathcal{P}$. We call a member of $\mathcal{C}$ a basic polar divisor of $\mathcal{P}$ if the corresponding Riemann form $E$ has the property that $E(D, D)=\mathbf{Z}$. Such a divisor $X$ is characterized by the property that every member of $\mathcal{C}$ is algebraically equivalent to $m X$ for some positive integer $m$, since two divisors on $A$ determine the same Riemann form if and only if they are algebraically equivalent (cf. also [S98, pp.27-28]).

Let $A^{\prime}$ be another abelian variety of dimension $n$, and $\left(\mathbf{C}^{n} / D^{\prime}, \xi^{\prime}\right)$ an analytic coordinate system of $A^{\prime}$. Every element $\lambda$ of $\operatorname{Hom}\left(A^{\prime}, A\right)$ coresponds to a $\mathbf{C}$-linear endomorphism $\Lambda$ of $\mathbf{C}^{n}$ such that $\Lambda D^{\prime} \subset D$ by the relation $\lambda \circ \xi^{\prime}=\xi \circ \Lambda$. Conversely every such $\Lambda$ determines an element $\lambda$ of $\operatorname{Hom}\left(A^{\prime}, A\right)$. Clearly $\lambda$ is an isogeny if and only if $\operatorname{det}(\Lambda) \neq 0$. Similarly every element of $\operatorname{Hom}_{\mathbf{Q}}\left(A^{\prime}, A\right)$ coresponds to an element of $\operatorname{End}\left(\mathbf{C}^{n}, \mathbf{C}\right)$ that sends $\mathbf{Q} D^{\prime}$ into $\mathbf{Q} D$, where $\mathbf{Q} D$ denotes the $\mathbf{Q}$-linear span of $D$.

Let $E$ be the Riemann form determined by a divisor $X$ on $A$, and $\lambda$ an isogeny of $A^{\prime}$ onto $A$. We can define a divisor $X^{\prime}$ on $A^{\prime}$ by $X^{\prime}=\lambda^{-1}(X)$. Then $X^{\prime}$ is also ample, and determines the Riemann form $E$ given by

$$
\begin{equation*}
E^{\prime}(z, w)=E(\Lambda z, \Lambda w) \tag{2.5}
\end{equation*}
$$

where $\Lambda$ is determined by $\lambda$ as above. This can be shown by considering $f \circ \Lambda$ with $f$ satisfying (2.4).

Let $\mathcal{P}=(A, \mathcal{C})$ and $\mathcal{P}^{\prime}=\left(A^{\prime}, \mathcal{C}^{\prime}\right)$ be two polarized abelian varieties of the same dimension. By an isogeny (resp. an isomorphism) of $\mathcal{P}^{\prime}$ onto $\mathcal{P}$ we understand an isogeny (resp. an isomorphism) $\lambda$ of $A^{\prime}$ onto $A$ such that $\lambda^{-1}(Y) \in \mathcal{C}^{\prime}$ for some $Y \in \mathcal{C}$. If that is so, then $\lambda^{-1}(X) \in \mathcal{C}^{\prime}$ for every $X \in \mathcal{C}$.
2.4. Whenever we consider an algebraic variety $V$, by our convention of $\S 2.1$, $V$ is defined in an affine or a projective space with a fixed coordinate system, so that we can speak of the coordinates of a point $x$ on $V$. For $x \in V$ and a field of definition $k$ for $V$ we denote by $k(x)$ the field generated over $k$ by the affine coordinates of $x$. (If $V$ is a projective variety, the affine coordinates of $x$ mean the quotients of the projective coordinates of $x$.)

We can also speak of the smallest field of definition for $V$. (See [W46, pp.71-72, Corollary 3 on page 71, in particular].) It is always finitely generated over the prime field. Also if $V$ is defined over a field $k$ and $\sigma$ is an isomorphism of $k$ onto a field $k^{\prime}$, Then $V^{\sigma}$ is well-defined; for example, it is defined by the equations which are the images under $\sigma$ of the $k$-rational defining equations for $V$. In particular, if $V$ is a point, $V^{\sigma}$ is the point whose coordinates are the images of the coordinates of $V$ under $\sigma$. If $f$ is a $k$-rational map of $V$ into an algebraic variety $V_{1}$ rational over $k$, then $f^{\sigma}$ is defined to be the rational map of $V^{\sigma}$ into $V_{1}^{\sigma}$ whose graph is the image of the graph of $f$ under $\sigma$. If $V$ is defined over $k$, we denote by $k(V)$ the field of all $k$-rational functions on $V$, that is, all $k$-rational maps of $V$ into the one-dimensional affine line. If $\sigma$ and $k^{\prime}$ are as above and $f \in k(V)$, then $f^{\sigma}$ is a well-defined element of $k^{\prime}\left(V^{\sigma}\right)$.

Let $\mathcal{P}=(A, \mathcal{C})$ be a polarized abelian variety. We say that $\mathcal{P}$ is defined (or rational) over $k$ and that $k$ is a field of definition (or rationality) for $\mathcal{P}$ if $A$ is defined over $k$ and $\mathcal{C}$ contains a $k$-rational divisor, say $X$. For $\sigma$ as above, we can define $\mathcal{P}^{\sigma}$ by $\mathcal{P}^{\sigma}=\left(A^{\sigma}, \mathcal{C}^{\sigma}\right)$, where $\mathcal{C}^{\sigma}$ is the polarization of $A^{\sigma}$ containing $X^{\sigma}$. (Here we need the fact that if $X$ is ample, so is $X^{\sigma}$.)

Suppose now $V$ is an algebraic variety defined over a subfield of C. Then, as we did in $\S 2.1$, we identify $V$ with the point set consisting of all the points with coordinates in $\mathbf{C}$ satisfying the defining equations for $V$. We denote by $\mathbf{C}(V)$ the union of $k(V)$ for all the subfields $k$ of $\mathbf{C}$. If $\sigma \in \operatorname{Aut}(\mathbf{C})$, then $V^{\sigma}$ as a point set consists of the images under $\sigma$ of all the points in $V$ in that sense. Let $k_{0}$ be the smallest field of definition for $V$. Then, for $\sigma, \tau \in \operatorname{Aut}(\mathbf{C})$ we have $V^{\sigma}=V^{\tau}$ if and only if $\sigma=\tau$ on $k_{0}$. Also, if $k$ is a subfield of $\mathbf{C}$ and $V^{\sigma}=V$ for every $\sigma \in \operatorname{Aut}(\mathbf{C} / k)$, then $V$ is defined over $k$.
2.5. Let $(A, \mathcal{C})$ and $\left(\mathbf{C}^{n} / D, \xi\right)$ be as above; let $E$ be the Riemann form on $\mathbf{C}^{n} / D$ determined by a divisor $X$ in $\mathcal{C}$. Given $\lambda \in \operatorname{End}_{\mathbf{Q}}(A)$, take $\Lambda \in \operatorname{End}\left(\mathbf{C}^{n}, \mathbf{C}\right)$ so that $\lambda \circ \xi=\xi \circ \Lambda$. Define an element $\Lambda^{\prime}$ of $\operatorname{End}\left(\mathbf{C}^{n}, \mathbf{C}\right)$ by

$$
\begin{equation*}
E\left(\Lambda^{\prime} x, y\right)=E(x, \Lambda y) \tag{2.6}
\end{equation*}
$$

From (2.1) we see that $\Lambda^{\prime}(\mathbf{Q} D) \subset \mathbf{Q} D$, and hence there is an element $\lambda^{\prime} \in \operatorname{End}_{\mathbf{Q}}(A)$ such that $\lambda^{\prime} \circ \xi=\xi \circ \Lambda^{\prime}$. We can easily verify that the map $\lambda \mapsto \lambda^{\prime}$ is a $\mathbf{Q}$-linear bijection of $\operatorname{End}_{\mathbf{Q}}(A)$ onto itself such that $\left(\lambda^{\prime}\right)^{\prime}=\lambda$ and $(\lambda \mu)^{\prime}=\mu^{\prime} \lambda^{\prime}$. We call this map the involution of $\operatorname{End}_{\mathbf{Q}}(A)$ determined by $X$, or by $\mathcal{C}$, since clearly it depends only on $\mathcal{C}$. Here are two easy facts:
(2.7) If $R$ (resp. $S$ ) is an element of $\mathbf{Q}_{2 n}^{2 n}$ (resp. $\mathbf{C}_{n}^{n}$ ) that represents $\Lambda$ with respect to a $\mathbf{Q}$-basis of $\mathbf{Q} D\left(\right.$ resp. C-basis of $\left.\mathbf{C}^{n}\right)$, then $R=T \cdot \operatorname{diag}[S, \bar{S}] T^{-1}$ for some $T \in G L_{2 n}(\mathbf{C})$ independent of $\Lambda$.

$$
\begin{equation*}
\operatorname{tr}\left(\Lambda^{\prime} \Lambda\right)>0 \text { if } \Lambda \neq 0 . \tag{2.8}
\end{equation*}
$$

Here $\operatorname{tr}(\Lambda)$ denotes the trace of $\Lambda$ as an $\mathbf{R}$-linear endomorphism. Indeed, since a $\mathbf{Q}$-basis of $\mathbf{Q} D$ gives an $\mathbf{R}$-basis of $\mathbf{C}^{n}$, we easily obtain (2.7). As for the latter, we note that $E\left(\Lambda^{\prime} x, i y\right)=E(x, i \Lambda y)$, which means that $\Lambda^{\prime}$ is the adjoint of $\Lambda$ with respect to the positive definite form of (2.3), so that we obtain (2.8).

Suppose now $A, X$, and $\lambda$ are rational over $k$; let $\sigma$ be an isomorphism of $k$ onto a field $k^{\prime}$. Let $\mu \mapsto \mu^{\prime \prime}$ be the involution of $\operatorname{End}_{\mathbf{Q}}\left(A^{\sigma}\right)$ determined by $X^{\sigma}$. Then
(2.9) $\lambda^{\prime}$ is rational over $k$ and $\left(\lambda^{\sigma}\right)^{\prime \prime}=\left(\lambda^{\prime}\right)^{\sigma}$.

This is because $\lambda^{\prime}$ can be defined algebro-geometrically by means of the Picard variety of $A$ without employing $E$, and we can let $\sigma$ act on the defining formula for $\lambda$. For details, see [S98, $\S 1.3$, formula (5) on page 5 in particular, and the last paragraph of $\S 3.3$ on page 25].
2.6. Lemma. Let $\left(\mathbf{C}^{n} / D, \xi\right)$ be an analytic coordinate system of an abelian variety $A$, and $X$ a divisor on $A$ corresponding to a Riemann form $E$ on $\mathbf{C}^{n} / D$ via $\xi$. For $s=\xi(x)$ and $t=\xi(y)$ with $x, y \in N^{-1} D, 0<N \in \mathbf{Z}$, put

$$
\begin{equation*}
\zeta_{X}(s, t)=\exp (2 \pi i N \cdot E(x, y)) \tag{2.10}
\end{equation*}
$$

Then $\zeta_{X}(s, t)$ is an $N$-th root of unity, and moreover, for every $\sigma \in \operatorname{Aut}(\mathbf{C})$,

$$
\begin{equation*}
\zeta_{X}(s, t)^{\sigma}=\zeta_{X^{\sigma}}\left(s^{\sigma}, t^{\sigma}\right) \tag{2.11}
\end{equation*}
$$

Furthermore, if $k$ is a field over which $A, X$, and all points on $A$ of finite order are rational, then $k$ contains the maximal abelian extension of $\mathbf{Q}$.

Proof. That $\zeta_{X}(s, t)$ is an $N$-th root of unity can be seen from the fact that $E(\Lambda, \Lambda) \subset \mathbf{Z}$. Formula (7) in $[\mathrm{S} 98, \mathrm{p} .24]$ shows that $\zeta_{X}(s, t)$ is the number $e_{X, N}(s, t)$ defined in $[S 98, \S 1.4]$ in a purely algebraic fashion without $E$. This definition is due to Weil [W48]. Therefore its behavior under $\sigma$ can easily be verified. Notice that if $E(\Lambda, \Lambda)=m \mathbf{Z}$ with $0<m \in \mathbf{Z}$ and $N$ is a multiple of $m$, then $\zeta_{X}(s, t)$ is a primitive $(N / m)$-th root of unity for suitable $s$ and $t$. Therefore we obtain the last assertion from (2.11).
2.7. We now generalize the above notion of polarized abelian variety by considering a structure $\mathcal{P}=\left(A, \mathcal{C}, \iota ;\left\{t_{i}\right\}_{i=1}^{r}\right)$ formed by a polarized abelian variety $(A, \mathcal{C})$ in the above sense, a ring-injection $\iota$ of a $\mathbf{Q}$-algebra $W$ (with identity element) into $\operatorname{End}_{\mathbf{Q}}(A)$, and an ordered set of points $\left\{t_{i}\right\}_{i=1}^{r}$ of $A$ of finite order. We always assume that $\iota(1)$ is the identity element of $\operatorname{End}(A)$. We say that $\mathcal{P}$ is defined (or rational) over a field $k$ and that $k$ is a field of definition (or rationality) for $\mathcal{P}$ if $(A, \mathcal{C})$, every element of $\iota(W) \cap \operatorname{End}(A)$, and every $t_{i}$ are all rational over $k$. We always take such a $k$ to be a subfield of $\mathbf{C}$. Given such a $k$ and an isomorphism $\sigma$ of $k$ onto a field (contained in $\mathbf{C}$ for the moment), we put $\mathcal{P}^{\sigma}=\left(A^{\sigma}, \mathcal{C}^{\sigma}, \iota^{\sigma} ;\left\{t_{i}^{\sigma}\right\}_{i=1}^{r}\right)$, where $\iota^{\sigma}(a)=\iota(a)^{\sigma}$. If $\mathcal{P}^{\prime}=\left(A^{\prime}, \mathcal{C}^{\prime}, \iota^{\prime} ;\left\{t_{i}^{\prime}\right\}_{i=1}^{r}\right)$ is another such structure with the same $W$, we understand by an isomorphism of $\mathcal{P}$ onto $\mathcal{P}^{\prime}$ an isomorphism $f$ of $(A, \mathcal{C})$ onto $\left(A^{\prime}, \mathcal{C}^{\prime}\right)$ such that $f \circ \iota(a)=\iota^{\prime}(a) \circ f$ for every $a \in W$ and $f\left(t_{i}\right)=t_{i}^{\prime}$ for every $i$. We call $f$ an automorphism of $\mathcal{P}$ if $\mathcal{P}=\mathcal{P}^{\prime}$, and denote by $\operatorname{Aut}(\mathcal{P})$ the group of all automorphisms of $\mathcal{P}$. Given an arbitrary $(A, \mathcal{C})$, we can construct $\mathcal{P}$ as above by taking $W=\mathbf{Q}$ and $\left\{t_{i}\right\}$ to be the set consisting of 0 . We shall alwlays identify such a $\mathcal{P}$ with $(A, \mathcal{C})$.
2.8. Theorem. (1) Given $\mathcal{P}$ as above, there exists a subfield $k_{1}$ of $\mathbf{C}$, called the field of moduli of $\mathcal{P}$, which is uniquely characterized by the following properties: (i) Every field of definition for $\mathcal{P}$ contains $k_{1}$; (ii) If $\mathcal{P}$ is defined over $k$ and $\sigma$ is an isomorphism of $k$ onto a subfield of $\mathbf{C}$, then $\sigma$ is the identity map on $k_{1}$ if and only if $\mathcal{P}^{\sigma}$ is isomorphic to $\mathcal{P}$.
(2) The field of moduli of $\mathcal{P}$ is algebraic over the field of moduli of $(A, \mathcal{C})$.
(3) If $\sum_{i=1}^{r} \mathbf{Z} t_{i} \supset\{t \in A \mid m t=0\}$ with an integer $m>2$, then $\mathcal{P}$ has a model rational over its field of moduli.
(4) $\mathcal{P}$ has a model over a finite algebraic extension of the field of moduli of $(A, \mathcal{C})$.

Assertion (1) is given in [S59] for $\mathcal{P}$ without $\left\{t_{i}\right\}$, and in [S65, §1.4] for $\mathcal{P}$ with $\left\{t_{i}\right\}$. Clearly the field of moduli of $\mathcal{P}$ is finitely generated over $\mathbf{Q}$. As for (2), see [S59, Proposition 8] and [S65, Proposition 1.11]; as for (3), see [S98, Proposition 21.1] and the remark after its proof, or [S65, Proposition 1.5]. Combining (2) and (3), we obtain (4).
2.9. The notation being as in $\S 2.7$, let $\psi$ be an equivalence class of $\mathbf{Q}$-linear representations of $W$ by complex matrices. We say that $(A, \iota)$ is of type $(W, \psi)$, if we can find an analytic coordinate system $\left(\mathbf{C}^{n} / D, \xi\right)$ of $A$ such that $\xi \circ \psi_{0}(a)=$ $\iota(a) \circ \xi$ for every $a \in W$ with a representation $\psi_{0}$ belonging to the class of $\psi$. In particular, take $W$ to be a CM-field $K$ in the sense of $\S 1.12$. If $A$ is of dimension $n$, then by (2.7) the direct sum of $\psi$ and its complex conjugate is equivalent to a rational representation of $K$ of degree $2 n$, which must be a multiple of a regular representation of $K$. Thus $[K: \mathbf{Q}]$ divides $2 n$. If $[K: \mathbf{Q}]=2 n$, then we see that $(K, \psi)$ is a CM-type in the sense of $\S 1.12$. Conversely, given a CM-type ( $K, \tau$ ), we can always find $(A, \iota)$ of type ( $K, \tau$ ); for details, see [S98, §6].

As an easy generalization we consider a $C M$-algebra $Y$, by which we understand a finite direct sum $Y=K_{1} \oplus \cdots \oplus K_{t}$ with CM-fields $K_{i}$. We take ( $A, \iota$ ) as above with $Y$ as $W$. Suppose that $2 \operatorname{dim}(A)=[Y: \mathbf{Q}]$; let $e_{i}$ be the identity element of $K_{i}$ and let $A_{i}=\iota\left(m_{i} e_{i}\right) A$ with a positive integer $m_{i}$ such that $\iota\left(m_{i} e_{i}\right) \in \operatorname{End}(A)$. We easily see that $A$ is isogenous to $A_{1} \times \cdots \times A_{t}$ and $\iota$ embeds $K_{i}$ into $\operatorname{End}_{\mathbf{Q}}\left(A_{i}\right)$. Denote this embedding by $\iota_{i}$. Since $\left[K_{i}: \mathbf{Q}\right] \leq 2 \operatorname{dim}\left(A_{i}\right)$ as observed above, and $[Y: \mathbf{Q}]=2 \operatorname{dim}(A)$, we obtain $\left[K_{i}: \mathbf{Q}\right]=2 \operatorname{dim}\left(A_{i}\right)$. Thus $\left(A_{i}, \iota_{i}\right)$ determines a CM-type ( $K_{i}, \Phi_{i}$ ). Then we see that ( $A, \iota$ ) is of type ( $Y, \Phi$ ) with $\Phi$ defined by

$$
\begin{equation*}
\Phi\left(\sum_{i=1}^{t} a_{i} e_{i}\right)=\operatorname{diag}\left[\Phi_{1}\left(a_{1}\right), \ldots, \Phi_{t}\left(a_{t}\right)\right] \quad\left(a_{i} \in K_{i}\right) \tag{2.12}
\end{equation*}
$$

## 3. Symmetric spaces and factors of automorphy

3.1. We first define our groups over $\mathbf{R}$, which will eventually be localizations of algebraic groups over a number field at archimedean primes. Using the notation of $\S 1.2$, we consider $S p(n, \mathbf{R}), G p(n, \mathbf{R})$ and the unitary groups of the following types:

$$
\begin{align*}
U\left(\eta_{n}\right) & =\left\{\alpha \in G L_{2 n}(\mathbf{C}) \mid \alpha^{*} \eta_{n} \alpha=\eta_{n}\right\},  \tag{3.1}\\
G U\left(\eta_{n}\right) & =\left\{\alpha \in G L_{2 n}(\mathbf{C}) \mid \alpha^{*} \eta_{n} \alpha=\nu(\alpha) \eta_{n} \text { with } \nu(\alpha) \in \mathbf{R}^{\times}\right\},  \tag{3.2}\\
U(m, n) & =\left\{\alpha \in G L_{r}(\mathbf{C}) \mid \alpha^{*} I_{m, n} \alpha=I_{m, n}\right\}, I_{m, n}=\operatorname{diag}\left[1_{m},-1_{n}\right],  \tag{3.3}\\
G U(m, n) & =\left\{\alpha \in G L_{r}(\mathbf{C}) \mid \alpha^{*} I_{m, n} \alpha=\nu(\alpha) I_{m, n} \text { with } \nu(\alpha) \in \mathbf{R}^{\times}\right\} . \tag{3.4}
\end{align*}
$$

Here $\eta_{n}$ is the matrix of (1.8), $m$ and $n$ are nonnegative integers and $r=$ $m+n>0$; we put $Z^{*}={ }^{t} \bar{Z}$ for a complex matrix $Z$. We easily see that if $\alpha$ belongs to $G p(n, A), G U\left(\eta_{n}\right)$ or $G U(m, n)$, then ${ }^{t} \alpha$ belongs to the same group and $\nu\left({ }^{t} \alpha\right)=\nu(\alpha)$. Clearly $U\left(\eta_{n}\right)$ and $G U\left(\eta_{n}\right)$ are isomorphic to $U(n, n)$ and $G U(n, n)$ respectively. For some technical reasons, however, we consider various objects in the unitary case in two different ways, with respect to these two types, even though the objects defined for the former types can easily be transferred to those defined for the latter types. Thus our discussion will be made in three cases, referred to as Case SP, Case UT, and Case UB: these correspond to $\operatorname{Sp}(n, \mathbf{R}), U\left(\eta_{n}\right)$, and $U(m, n)$.

We now define domains $\mathfrak{H}_{n}, \mathcal{H}_{n}$, and $\mathfrak{B}_{m . n}$ by

$$
\begin{align*}
& \mathfrak{H}_{n}=\left\{\left.z \in \mathbf{C}_{n}^{n}\right|^{t} z=z, \operatorname{Im}(z)>0\right\}  \tag{3.5}\\
& \mathcal{H}_{n}=\left\{z \in \mathbf{C}_{n}^{n} \mid i\left(z^{*}-z\right)>0\right\}  \tag{3.6}\\
& \mathfrak{B}_{m, n}=\mathfrak{B}(m, n)=\left\{z \in \mathbf{C}_{n}^{m} \mid 1_{n}-z^{*} z>0\right\} \quad(m n>0) \tag{3.7}
\end{align*}
$$

Here for a hermitian matrix $\xi$ (in particular, for a real symmetric matrix $\xi$ ) we write $\xi>0$ (resp. $\xi \geq 0$ ) to indicate that $\xi$ is positive definite (resp. nonnegative). We then put

$$
\begin{gather*}
S^{n}=\left\{h \in \mathbf{C}_{n}^{n} \mid h^{*}=h\right\}, \quad S_{n}^{+}=\left\{h \in S^{n} \mid h>0\right\},  \tag{3.8}\\
B(z)=\left[\begin{array}{cc}
z^{*} & z \\
1_{n} & 1_{n}
\end{array}\right], \quad \xi(z)=i\left(\bar{z}-{ }^{t} z\right), \quad \eta(z)=i\left(z^{*}-z\right)  \tag{3.9}\\
\left(z \in \mathbf{C}_{n}^{n}, \text { Cases SP, UT }\right), \\
B(z)=\left[\begin{array}{cc}
1_{m} & z \\
z^{*} & 1_{n}
\end{array}\right], \quad \xi(z)=1_{m}-\bar{z} \cdot{ }^{t} z, \quad \eta(z)=1_{n}-z^{*} z  \tag{3.10}\\
\left(z \in \mathbf{C}_{n}^{m}, \text { Case UB }\right) .
\end{gather*}
$$

For the moment we assume $m n \neq 0$. We shall discuss the case $m n=0$ later. By straightforward calculations we can verify that

$$
\begin{array}{ll}
i \cdot B(z)^{*} \eta_{n} B(z)=\operatorname{diag}[t \xi(z),-\eta(z)] & \text { (Cases SP, UT), } \\
B(z)^{*} I_{m, n} B(z)=\operatorname{diag}\left[{ }^{t} \xi(z),-\eta(z)\right] & \text { (Case UB) } \\
\operatorname{det}[B(z)]= \begin{cases}\operatorname{det}\left(z^{*}-z\right) & \text { (Cases SP, UT), } \\
\operatorname{det}[\xi(z)]=\operatorname{det}[\eta(z)] & \text { (Case UB). }\end{cases} \tag{3.13}
\end{array}
$$

Now if $\eta(z)>0$ in Case UB, then $\operatorname{det}[B(z)] \neq 0$, and hence the left-hand side of (3.12) has signature $(m, n)$, so that $\xi(z)>0$. Similarly $\eta(z)>0$ if $\xi(z)>0$.
3.2. Lemma. Case UT: Let $\mathfrak{X}_{u}$ be the set of all matrices $X$ in $\mathbf{C}_{2 n}^{2 n}$ such that $i X^{*} \eta_{n} X=\operatorname{diag}[v,-w]$ with $v, w \in S_{n}^{+}$. Then the map $(z, \lambda, \mu) \mapsto B(z) \operatorname{diag}[\bar{\lambda}, \mu]$ gives a bijection of $\mathcal{H}_{n} \times G L_{n}(\mathbf{C}) \times G L_{n}(\mathbf{C})$ onto $\mathfrak{X}_{u}$.

Case SP: Let $\mathfrak{X}$ be the set of all $X$, belonging to the set $\mathfrak{X}_{u}$ defined above, of the form $X=\left[\begin{array}{ll}\bar{y} & y\end{array}\right]$ with $y \in \mathbf{C}_{n}^{2 n}$. Then the map $(z, \mu) \mapsto B(z) \operatorname{diag}[\bar{\mu}, \mu]$ gives a bijection of $\mathfrak{H}_{n} \times G L_{n}(\mathbf{C})$ onto $\mathfrak{X}$.

Case UB: Let $\mathfrak{X}$ be the set of all $X \in \mathbf{C}_{m+n}^{m+n}$ such that $X^{*} I_{m, n} X=\operatorname{diag}[v,-w]$ with $v \in S_{m}^{+}$and $w \in S_{n}^{+}$. Then the map $(z, \lambda, \mu) \mapsto B(z) \operatorname{diag}[\bar{\lambda}, \mu]$ gives a bijection of $\mathfrak{B}_{m, n} \times G L_{m}(\mathbf{C}) \times G L_{n}(\mathbf{C})$ onto $\mathfrak{X}$.

This lemma is completely elementary. The proof in Case SP was given in [S97, Lemma 7.10]. Cases UT nd UB can be proved similarly.
3.3. Write simply $\mathfrak{X}$ for the set $\mathfrak{X}_{u}$ defined in Case UT in the above lemma. Thus we have $\mathfrak{X}$ defined in each of the three cases SP, UT, and UB. If $\alpha \in G p(n, \mathbf{R})$ (resp. $\alpha \in G U\left(\eta_{n}\right), \alpha \in G U(m, n), m n>0$ ) with $\nu(\alpha)>0$, then we easily see that $\alpha \mathfrak{X} \subset \mathfrak{X}$. Given such an $\alpha$, we define the action of $\alpha$ on $\mathfrak{H}_{n}$ (resp. $\mathcal{H}_{n}, \mathfrak{B}_{m, n}$ ) and two factors of automorphy $\lambda(\alpha, z)$ and $\mu(\alpha, z)$ as follows. First we consider Case UB. Let $\alpha \in G U(m, n)$ and $z \in \mathfrak{B}_{m, n}$. By the above lemma $B(z) \in \mathfrak{X}$, and so $\alpha B(z) \in \mathfrak{X}$. By the same lemma we can put $\alpha B(z)=B(w) \operatorname{diag}[\bar{\lambda}, \mu]$ with unique $w \in \mathfrak{B}_{m, n}, \lambda \in G L_{m}(\mathbf{C})$, and $\mu \in G L_{n}(\mathbf{C})$. Let us put $w=\alpha z=\alpha(z)$, $\lambda=\lambda(\alpha, z)$, and $\mu=\mu(\alpha, z)$. Then we have

$$
\begin{equation*}
\alpha B(z)=B(\alpha z) \operatorname{diag}[\overline{\lambda(\alpha, z)}, \mu(\alpha, z)] \tag{3.14}
\end{equation*}
$$

We can do the same in the other two cases. In Case UT both $\lambda(\alpha, z)$ and $\mu(\alpha, z)$ belong to $G L_{n}(\mathbf{C})$; in Case SP we have $\lambda(\alpha, z)=\mu(\alpha, z) \in G L_{n}(\mathbf{C})$. To make our formulas short, we shall often put $\lambda(\alpha, z)=\lambda_{\alpha}(z)$ and $\mu(\alpha, z)=\mu_{\alpha}(z)$. For $\alpha=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right]$ with $d$ of size $n$ let us put $a=a_{\alpha}, b=b_{\alpha}, c=c_{\alpha}$, and $d=d_{\alpha}$. Then from (3.14) we immediately see that $a_{\alpha} z+b_{\alpha}=\alpha(z) \mu_{\alpha}(z)$ and $\mu_{\alpha}(z)=c_{\alpha} z+d_{\alpha}$. Since $\mu$ is invertible, we obtain

$$
\begin{array}{cc}
\alpha z=\left(a_{\alpha} z+b_{\alpha}\right)\left(c_{\alpha} z+d_{\alpha}\right)^{-1} & \text { (all cases) } \\
\lambda_{\alpha}(z)=\bar{c}_{\alpha} \cdot{ }^{t} z+\bar{d}_{\alpha}, \quad \mu_{\alpha}(z)=c_{\alpha} z+d_{\alpha} & \text { (Cases SP, UT), } \\
\lambda_{\alpha}(z)=\bar{b}_{\alpha} \cdot{ }^{t} z+\bar{a}_{\alpha}, \quad \mu_{\alpha}(z)=c_{\alpha} z+d_{\alpha} & \text { (Case UB). } \tag{3.17}
\end{array}
$$

The first formula means that we can let $\alpha$ act on $\mathfrak{H}_{n}, \mathcal{H}_{n}$, or $\mathfrak{B}_{m, n}$ by defining $\alpha z$ by (3.15). Notice that in all three cases $\lambda_{\alpha}(z)$ and $\mu_{\alpha}(z)$ are holomorphic in $z$. Applying another element $\beta$ with $\nu(\beta)>0$ to (3.14), we see that $\beta(\alpha z)=(\beta \alpha) z$ and

$$
\begin{equation*}
\lambda(\beta \alpha, z)=\lambda(\beta, \alpha z) \lambda(\alpha, z), \quad \mu(\beta \alpha, z)=\mu(\beta, \alpha z) \mu(\alpha, z) \tag{3.18}
\end{equation*}
$$

From (3.11), (3.12), and (3.14) we obtain

$$
\begin{equation*}
\lambda_{\alpha}(z)^{*} \xi(\alpha z) \lambda_{\alpha}(z)=\nu(\alpha) \xi(z), \quad \mu_{\alpha}(z)^{*} \eta(\alpha z) \mu_{\alpha}(z)=\nu(\alpha) \eta(z) \tag{3.19}
\end{equation*}
$$

Define a scalar factor of automorphy $j_{\alpha}(z)$ and a real-valued function $\delta$ by

$$
\begin{align*}
j_{\alpha}(z) & =j(\alpha, z)=\operatorname{det}\left[\mu_{\alpha}(z)\right]  \tag{3.20}\\
\delta(z) & =\left\{\begin{array}{lc}
\operatorname{det}\left[2^{-1} \eta(z)\right] & \text { (all cases) } \\
\operatorname{det}[\eta(z)] & \text { (Case UB) }
\end{array}\right. \tag{3.21}
\end{align*}
$$

From (3.18) we obtain $j_{\beta \alpha}(z)=j_{\beta}(\alpha z) j_{\alpha}(z)$ and

$$
\begin{equation*}
\delta(\alpha z)=\nu(\alpha)^{n}\left|j_{\alpha}(z)\right|^{-2} \delta(z) \tag{3.22}
\end{equation*}
$$

We note here an easy but essential relation:

$$
\begin{equation*}
\operatorname{det}\left[\lambda_{\alpha}(z)\right]=\operatorname{det}(\bar{\alpha}) \nu(\alpha)^{-n} \cdot j_{\alpha}(z) \quad(\text { Cases UT, UB), } \tag{3.23}
\end{equation*}
$$

which follows immediately from (3.13), (3.14), and (3.19).
So far we have assumed $m n>0$ in Case UB. We now make the following convention: if $m n=0$ in Case UB , then $\mathfrak{B}(m, n)$ consists of the single element 0 , our group acts on it trivially, and $B(0)=1_{m+n} ; G L_{0}(\mathbf{C})$ denotes the group consisting only of the identity element 1 , and $\operatorname{det}(1)=1 ;\left(\lambda_{\alpha}(z), \mu_{\alpha}(z), j_{\alpha}(z)\right)$ and $(\xi(z), \eta(z), \delta(z))$ are elements of $G L_{m}(\mathbf{C}) \times G L_{n}(\mathbf{C}) \times \mathbf{C}^{\times}$determined by

$$
\begin{gather*}
\xi(z)=1_{m}, \quad \lambda_{\alpha}(z)=\bar{\alpha}, \quad \text { and } \quad j_{\alpha}(z)=1 \quad \text { if } \quad n=0  \tag{3.24a}\\
\eta(z)=1_{n}, \quad \mu_{\alpha}(z)=\alpha, \quad \text { and } \quad j_{\alpha}(z)=\operatorname{det}(\alpha) \quad \text { if } m=0,  \tag{3.24b}\\
\delta(z)=1 \quad \text { if } m n=0 .
\end{gather*}
$$

Thus $\eta(z)=\mu_{\alpha}(z)=1$ if $n=0$, and $\xi(z)=\lambda_{\alpha}(z)=1$ if $m=0$.
Also we denote by $\mathbf{C}^{0}$ the 0 -dimensional vector space $\{0\}$ and let $G L_{0}(\mathbf{C})$ act on $\mathbf{C}^{0}$ trivially. We identify $\mathbf{C}^{m} \times \mathbf{C}^{n}$ and $G L_{m}(\mathbf{C}) \times G L_{n}(\mathbf{C})$ with $\mathbf{C}^{m+n}$ and $G L_{m+n}(\mathbf{C})$ in an obvious way, if $m n=0$. Then we denote by $\operatorname{diag}[a, b]$ the element of $G L_{m+n}(\mathbf{C})$ identified with $(a, b) \in G L_{m}(\mathbf{C}) \times G L_{n}(\mathbf{C})$. Then (3.14), (3.18), (3.19), (3.22), and (3.23) are valid.
3.4. Lemma. (1) Let $d z=\left(d z_{h k}\right)$ be a matrix of the same shape as $z$ whose entries are 1 -forms $d z_{h k}$ on each space. Then for $\alpha$ in $S P(n, \mathbf{R}), U\left(\eta_{n}\right)$, or in $U(q, r)$ we have $d(\alpha z)={ }^{t} \lambda_{\alpha}(z)^{-1} \cdot d z \cdot \mu_{\alpha}(z)^{-1}$.
(2) The jacobian of the map $z \mapsto \alpha z$ for such an $\alpha$ is $j_{\alpha}(z)^{-\kappa}$, where $\kappa=$ $n+1,2 n$, and $m+n$ in Cases $S P, U T$, and $U B$, respectively.
(3) Define a differential form $\mathbf{d} z$ on our space by

$$
\mathbf{d} z= \begin{cases}\delta(z)^{-n-1} \prod_{h \leq k}\left[(i / 2) d z_{h k} \wedge d \bar{z}_{h k}\right] & (\text { Case SP }) \\ \delta(z)^{-m-n} \prod_{h=1}^{m} \prod_{k=1}^{n}\left[(i / 2) d z_{h k} \wedge d \bar{z}_{h k}\right] & (\text { Cases UT and UB) }\end{cases}
$$

where we put $m=n$ in Case UT. Then $\mathbf{d} z$ is invariant under the map of (2), and therefore it defines an invariant measure on our space.

Proof. We first consider Case UB. For $z, w \in \mathfrak{B}_{m, n}$ we have

$$
B(w)^{*} I_{m, n} B(z)=\left[\begin{array}{cc}
1-w z^{*} & z-w \\
w^{*}-z^{*} & w^{*} z-1
\end{array}\right]
$$

Changing $w$ and $z$ for $\alpha w$ and $\alpha z$ and employing (3.14), we find that

$$
\begin{equation*}
\alpha z-\alpha w={ }^{t} \lambda_{\alpha}(w)^{-1}(z-w) \mu_{\alpha}(z)^{-1} \tag{3.25a}
\end{equation*}
$$

Taking similarly $B(w)^{*} \eta_{n} B(z)$ in Cases SP and UT, we see that (3.25a) is true in all cases. From this we obtain (1). Computing the determinant of the linear map $x \mapsto{ }^{t} \lambda_{\alpha}(z)^{-1} x \mu_{\alpha}(z)^{-1}$, we obtain (2), which together with (3.22) proves (3). Clearly (3.25a) combined with (3.23) implies

$$
\begin{equation*}
\operatorname{det}(\alpha z-\alpha w)=\operatorname{det}(\alpha) j_{\alpha}(z)^{-1} j_{\alpha}(w)^{-1} \operatorname{det}(z-w) \tag{3.25b}
\end{equation*}
$$

3.5. We now take a totally real algebraic number field $F$ of finite degree and denote by a the set of all archimedean primes of $F$. In Cases UT and UB we take a CM-type ( $K, \tau$ ) as in $\S 1.12$ with $K$ containing $F$ as its maximal real subfield, and denote by $\rho$ the generator of $\operatorname{Gal}(K / F)$. We shall also employ the symbols $\mathfrak{g}, \mathfrak{r}, \mathbf{h}$, and $\mathbf{k}$ introduced in $\S 1.4$ and $\S 1.8$. Then $\tau$ can be written $\tau=\left\{\tau_{v}\right\}_{v \in \mathbf{a}}$ with an embedding $\tau_{v}: K \rightarrow \mathbf{C}$ which coincides with $v$ on $F$. Hereafter we fix $\tau$ and for $a \in K$ denote by $a_{v}$ the image of $a$ under $\tau_{v}$. Then we identify a with $\tau$ and view a also as the set of all archimedean primes of $K$. For $c \in F_{\mathbf{A}}^{\times}$we write $c \gg 0$ if $c_{v}>0$ for every $v \in \mathbf{a}$. To make our exposition uniform, we make a convention that $K=F$ and $\mathfrak{r}=\mathfrak{g}$ in Case SP; thus we use $K$ and $\mathfrak{r}$ in all three cases.

Given a set $X$, we denote by $X^{\mathbf{a}}$ the product of a copies of $X$, that is, the set of all indexed elements $\left(x_{v}\right)_{v \in \mathbf{a}}$ with $x_{v}$ in $X$. Then all the embeddings of $F$ into $\mathbf{R}$ (resp. $K$ into $\mathbf{C}$ ) given by the elements of a determine an isomorphism of $F \otimes_{\mathbf{Q}} \mathbf{R}$ onto $\mathbf{R}^{\mathbf{a}}$ (resp. $K \otimes_{\mathbf{Q}} \mathbf{R}$ onto $\mathbf{C}^{\mathbf{a}}$.). Similarly we obtain embeddings of $F_{n}^{m}$ and $K_{n}^{m}$ into $\left(\mathbf{R}_{n}^{m}\right)^{\mathbf{a}}$ and $\left(\mathbf{C}_{n}^{m}\right)^{\mathbf{a}}$. We view the former sets as subsets of the latter sets. Thus for $\alpha \in F_{n}^{m}$ and $v \in \mathbf{a}$ the $v$-component $\alpha_{v}$ of $\alpha$ considered in $\left(\mathbf{R}_{n}^{m}\right)^{\mathbf{a}}$ is the $v$-th conjugate of $\alpha$. If $\alpha \in K_{n}^{m}$, then $\alpha_{v}$ is the image of $\alpha$ under $\tau_{v}$.

We now define algebraic groups $\widetilde{G}, G$, and $G_{0}$ by

$$
\begin{array}{lll}
\widetilde{G}=G p(n, F), & G=S p(n, F) & (\text { Case SP), } \\
\widetilde{G}=G U\left(\eta_{n}\right), & G=U\left(\eta_{n}\right) & (\text { Case UT) } \tag{3.27}
\end{array}
$$

$$
\begin{array}{cl}
\widetilde{G}=G U(\mathcal{T}), \quad G=U(\mathcal{T}) & (\text { Case UB }), \\
G_{0}=\left\{\alpha \in \widetilde{G} \mid \operatorname{det}(\alpha)=\nu(\alpha)^{n}, \nu(\alpha) \in \mathbf{Q}\right\} & (\text { Cases SP and UT) } \tag{3.29}
\end{array}
$$

Here we are using the notation of (1.6), (1.7), (1.8), (1.9), and (1.10) with the present $(F, K, \rho) ; \mathcal{T}$ is a fixed element of $G L_{r}(K)$ such that $\mathcal{T}^{*}=-\mathcal{T}$; for a matrix $Z$ we put $Z^{*}={ }^{t} Z^{\rho}$ as we did in (1.5). Recall that in each case we have a homomorphism $\nu: \widetilde{G} \rightarrow F^{\times}$. Though Case UT is included in Case UB, we treat them in different ways. In Case SP we can ignore the condition $\operatorname{det}(\alpha)=\nu(\alpha)^{n}$ in the definition of $G_{0}$, since it is true for every $\alpha \in \widetilde{G}$ as noted in (1.11).

For each $v \in \mathbf{v}$ we have localizations $\widetilde{G}_{v}$ and $G_{v}$ of $\widetilde{G}$ and $G$; recall that $\widetilde{G}_{\mathbf{a}}=$ $\prod_{v \in \mathbf{a}} \widetilde{G}_{v}$ and $G_{\mathbf{a}}=\prod_{v \in \mathbf{a}} G_{v}$. In Case UB let $\left(m_{v}, n_{v}\right)$ be the signature of $i \mathcal{I}_{v}$. We then put

$$
\begin{gather*}
\widetilde{G}_{\mathbf{A}+}=\left\{\alpha \in G_{\mathbf{A}} \mid \nu(\alpha)_{\mathbf{a}} \in F_{\mathbf{a}+}^{\times}\right\}, \quad F_{\mathbf{a}+}^{\times}=\left\{x \in F_{\mathbf{a}}^{\times} \mid x \gg 0\right\},  \tag{3.30}\\
\widetilde{G}_{\mathbf{a}+}=\widetilde{G}_{\mathbf{a}} \cap \widetilde{G}_{\mathbf{A}+}, \quad \widetilde{G}_{+}=\widetilde{G} \cap \widetilde{G}_{\mathbf{A}+},  \tag{3.31a}\\
\left(G_{0}\right)_{\mathbf{A}+}=\left(G_{0}\right)_{\mathbf{A}} \cap \widetilde{G}_{\mathbf{A}+}, \quad G_{0+}=G_{0} \cap \widetilde{G}_{\mathbf{A}+} . \tag{3.31b}
\end{gather*}
$$

In particular, in Case UB we put

$$
\begin{equation*}
\widetilde{\mathfrak{G}}=\prod_{v \in \mathbf{a}} G U\left(m_{v}, n_{v}\right), \quad \mathfrak{G}=\prod_{v \in \mathbf{a}} U\left(m_{v}, n_{v}\right), \quad \widetilde{\mathfrak{G}}_{+}=\{\alpha \in \widetilde{\mathfrak{G}} \mid \nu(\alpha) \gg 0\} \tag{3.32}
\end{equation*}
$$

We define a space $\mathcal{H}$ by

$$
\mathcal{H}= \begin{cases}\left(\mathfrak{H}_{n}\right)^{\mathbf{a}} & (\text { Case SP) }  \tag{3.33}\\ \left(\mathcal{H}_{n}\right)^{\mathbf{a}} & (\text { Case UT) } \\ \prod_{v \in \mathbf{a}} \mathfrak{B}\left(m_{v}, n_{v}\right) & (\text { Case UB })\end{cases}
$$

We now define the action of the elements of $\widetilde{G}_{\mathbf{A}+}$ on $\mathcal{H}$ as follows: In Cases SP and UT, given $\xi=\left(\xi_{v}\right)_{v \in \mathbf{v}} \in \widetilde{G}_{\mathbf{A}+}$ and $z=\left(z_{v}\right)_{v \in \mathbf{a}} \in \mathcal{H}$, we consider $\xi_{\mathbf{a}}=\left(\xi_{v}\right)_{v \in \mathbf{a}}$ and put $\xi(z)=\left(\xi_{v}\left(z_{v}\right)\right)_{v \in \mathbf{a}}$. In Case UB we take an element $Q_{v} \in G L_{m}(\overline{\mathbf{Q}})$ so that

$$
\begin{equation*}
i \mathcal{T}_{v}=Q_{v} I_{m_{v}, n_{v}} Q_{v}^{*} \tag{3.34}
\end{equation*}
$$

(The reason why we take $G L_{m}(\overline{\mathbf{Q}})$ instead of $G L_{m}(\mathbf{C})$ will be explained later.) Clearly the map

$$
\begin{equation*}
\alpha \mapsto\left(Q_{v}^{-1} \alpha_{v} Q_{v}\right)_{v \in \mathbf{a}} \tag{3.35}
\end{equation*}
$$

sends $\widetilde{G}_{\mathbf{A}}$ into $\tilde{\mathfrak{G}}$. We then let $\widetilde{G}_{\mathbf{A}+}$ act on $\mathcal{H}$ through the map $\widetilde{G}_{\mathbf{A}+} \rightarrow \widetilde{\mathfrak{G}}_{+}$. For $\alpha \in \widetilde{G}_{\mathbf{A}+}$ and $z=\left(z_{v}\right)_{v \in \mathbf{a}} \in \mathcal{H}$ we put

$$
\begin{align*}
& \lambda_{v}(\alpha, z)=\lambda\left(\alpha_{v}, z_{v}\right), \quad \mu_{v}(\alpha, z)=\mu\left(\alpha_{v}, z_{v}\right) \quad \text { (Cases SP, UT) }  \tag{3.36}\\
& \lambda_{v}(\alpha, z)=\lambda\left(Q_{v}^{-1} \alpha_{v} Q_{v}, z_{v}\right), \quad \mu_{v}(\alpha, z)=\mu\left(Q_{v}^{-1} \alpha_{v} Q_{v}, z_{v}\right) \quad \text { (Case UB). } \tag{3.37}
\end{align*}
$$

Then, in Case UB, from (3.14) we obtain

$$
\begin{equation*}
\alpha_{v} Q_{v} B\left(z_{v}\right)=Q_{v} B\left((\alpha z)_{v}\right) \cdot \operatorname{diag}\left[\overline{\lambda_{v}(\alpha, z)}, \mu_{v}(\alpha, z)\right] \tag{3.38}
\end{equation*}
$$

3.6. It is sometimes necessary to send $\mathfrak{H}_{n}$ and $\mathcal{H}_{n}$ onto bounded domains. In Case UT the domain $\mathcal{H}_{n}$ can be sent onto $\mathfrak{B}(n, n)$ as will be shown below. To deal with problems of this type and to define the domain in Case SP, we put

$$
\begin{align*}
& \mathfrak{B}_{n}=\left\{\left.z \in \mathfrak{B}(n, n)\right|^{t} z=z\right\}  \tag{3.39}\\
& \mathfrak{E}=\left\{\alpha \in G L_{2 n}(\mathbf{C}) \mid \alpha^{*} \eta_{n} \alpha=-i I_{n, n}\right\},  \tag{3.40}\\
& G^{\prime}=S p(n, \mathbf{C}) \cap U(n, n), \quad \mathfrak{E}^{\prime}=S p(n, \mathbf{C}) \cap \mathfrak{E} . \tag{3.41}
\end{align*}
$$

Clearly $\mathfrak{E}=U\left(\eta_{n}\right) \alpha=\alpha U(n, n)$ for every $\alpha \in \mathfrak{E}$ and $\mathfrak{E}^{\prime}=S p(n, \mathbf{R}) \beta=\beta G^{\prime}$ for every $\beta \in \mathfrak{E}^{\prime}$. As an exemplary element of $\mathfrak{E}^{\prime}$ we note

$$
\beta_{0}=\left[\begin{array}{cc}
\varepsilon 1_{n} & \varepsilon 1_{n}  \tag{3.42}\\
-\bar{\varepsilon} 1_{n}^{\prime} & \bar{\varepsilon} 1_{n}
\end{array}\right]
$$

with a complex number $\varepsilon$ such that $\varepsilon^{2}=i / 2$.
Now, given $\alpha \in \mathfrak{E}$ and $z \in \mathfrak{B}_{n, n}$, we see that $\alpha B(z)$ belongs to the set $\mathfrak{X}_{u}$ of Lemma 3.2 in Case UT, and so, by that lemma, $\alpha B(z)=\left[\begin{array}{cc}w^{*} & w \\ 1 & 1\end{array}\right] \operatorname{diag}[\lambda, \mu]$ with $\lambda, \mu \in G L_{n}(\mathbf{C})$ and $w \in \mathcal{H}_{n}$. Then we put $w=\alpha z, \lambda_{\alpha}(z)=i \bar{\lambda}$, and $\mu_{\alpha}(z)=\mu$. Thus

$$
\alpha\left[\begin{array}{cc}
1_{n} & z  \tag{3.43}\\
z^{*} & 1_{n}
\end{array}\right]=\left[\begin{array}{cc}
(\alpha z)^{*} & \alpha z \\
1_{n} & 1_{n}
\end{array}\right]\left[\begin{array}{cc}
i \cdot \overline{\lambda_{\alpha}(z)} & 0 \\
0 & \mu_{\alpha}(z)
\end{array}\right] \quad\left(z \in \mathfrak{B}_{n, n}\right) .
$$

In this way $z \mapsto \alpha z$ sends $\mathfrak{B}_{n, n}$ into $\mathcal{H}_{n}$. A simple calculation shows that $\alpha z=$ $\left(a_{\alpha} z+b_{\alpha}\right)\left(c_{\alpha} z+d_{\alpha}\right)^{-1}$.

Since $\varepsilon=i \bar{\varepsilon}$ and $\mathfrak{E}^{\prime}=S p(n, \mathbf{R}) \beta_{0}$ with $\beta_{0}$ of (3.42), we easily see that every element of $\mathfrak{E}^{\prime}$ is of the form $[i \bar{p} p]$ with $p \in \mathbf{C}_{n}^{2 n}$. Therefore if $\alpha \in \mathcal{E}^{\prime}$ and ${ }^{t} z=z$, then we see that $\alpha B(z)=\left[\begin{array}{ll}i \bar{y} & y\end{array}\right]$ with $y \in \mathbf{C}_{n}^{2 n}$. From this and Lemma 3.2 (Case SP ) we see that $\lambda_{\alpha}(z)=\mu_{\alpha}(z)$ and $\alpha z \in \mathfrak{H}_{n}$. Thus $\alpha$ sends $\mathfrak{B}_{n}$ into $\mathfrak{H}_{n}$.

Considering similarly $\alpha^{-1} B(w)$ with $w \in \mathcal{H}_{n}$, we find that every $\alpha$ in $\mathfrak{E}$ gives a bijection of $\mathfrak{B}_{n, n}$ onto $\mathcal{H}_{n}$, and also a bijection of $\mathfrak{B}_{n}$ onto $\mathfrak{H}_{n}$ if $\alpha \in \mathfrak{E}^{\prime}$. The action is associative in the sense that $(\alpha \gamma) z=\alpha(\gamma z)$ for $\gamma \in U(n, n)$ and $(\beta \alpha) z=\beta(\alpha z)$ for $\beta \in U\left(\eta_{n}\right)$. Furthermorer, (3.18) holds for $(\beta, \alpha) \in \mathfrak{E} \times U(n, n)$ and also for $(\beta, \alpha) \in U\left(\eta_{n}\right) \times \mathfrak{E} ;(3.19)$ is valid for $\alpha \in \mathfrak{E}$ and $z \in \mathfrak{B}_{n, n}$ if we understand that $\xi(\alpha z)={ }^{t} \eta(\alpha z)=i\left(\overline{\alpha z}-{ }^{t}(\alpha z)\right)$, as defined in (3.9). Finally we can show that (3.25a) is true for $z, w \in \mathfrak{B}_{n, n}$ and $\alpha \in \mathfrak{E}^{\prime}$, and hence

$$
\begin{equation*}
d(\alpha z)={ }^{t} \lambda_{\alpha}(z)^{-1} \cdot d z \cdot \mu_{\alpha}(z)^{-1} \quad\left(\alpha \in \mathfrak{E}, z \in \mathfrak{B}_{n, n}\right) . \tag{3.44}
\end{equation*}
$$

These are essentially special cases of [S97, Lemma A2.3] and the formulas stated in its proof.

## 4. Families of polarized abelian varieties

4.1. Let us now consider a structure $\mathcal{P}=\left(A, \mathcal{C}, \iota ;\left\{t_{i}\right\}_{i=1}^{s}\right)$ of $\S 2.7$ under the following conditions:
(4.1) $\iota$ is a ring-injection of a field $K$ into $\operatorname{End}_{\mathbf{Q}}(A)$.
(4.2) $\iota(K)$ is stable under the involution $\alpha \mapsto \alpha^{\prime}$ of $\operatorname{End}_{\mathbf{Q}}(A)$ determined by $\mathcal{C}$ as in §2.5.

Here $K$ is as in $\S 3.5$; namely $K=F$ in Case SP, and $K$ is a CM-field in Cases UT and UB. Given $\mathcal{P}$ satisfying these conditions, let $d$ be the dimension of $A$. Let $\left(\mathbf{C}^{d} / \Lambda, \xi\right)$ be an analytic coordinate system for $A$ in the sense of $\S 2.2$. Then we find a ring-injection $\Psi: K \rightarrow \mathbf{C}_{d}^{d}$ given by $\iota(a) \xi(u)=\xi(\Psi(a) u)$ for $a \in K$ and $u \in \mathbf{C}^{d}$. Let $\mathbf{Q} \Lambda$ denote the $\mathbf{Q}$-linear span of $\Lambda$ in $\mathbf{C}^{d}$. Then $\mathbf{Q} \Lambda$ is a $2 d$-dimensional vector space over $\mathbf{Q}$, and is stable under $\Psi(K)$. Thus it has a structure of a vector space
over $K$, so that $\mathbf{Q} \Lambda$ is isomorphic to $K_{r}^{1}$ with the integer $r$ such that $2 d=r[K: \mathbf{Q}]$. Since $K_{\mathbf{a}}=K \otimes_{\mathbf{Q}} \mathbf{R}$, the isomorphism of $K_{r}^{1}$ onto $\mathbf{Q} \Lambda$ can be extended to an $\mathbf{R}$ linear isomorphism $q$ of $\left(K_{\mathbf{a}}\right)_{r}^{1}$ onto $\mathbf{C}^{d}$ such that $q(a x)=\Psi(a) q(x)$ for $a \in K$ and $x \in K_{r}^{1}$. Let $L$ be the inverse image of $\Lambda$. Then we obtain a commutative diagram:


Let $E_{X}(x, y)$ be the Riemann form determined by a divisor $X$ in $\mathcal{C}$, and $\rho$ the involution of $K$ such that $\iota\left(a^{\rho}\right)=\iota(a)^{\prime}$, where $\alpha^{\prime}$ is as in (4.2). Then, by (2.6),

$$
\begin{equation*}
E_{X}(\Psi(a) u, v)=E_{X}\left(u, \Psi\left(a^{\rho}\right) v\right) \quad \text { for every } \quad a \in W \tag{4.4}
\end{equation*}
$$

Let $\operatorname{tr}(\Psi(a))$ denote the trace of $\Psi(a)$ as an R-linear map. Clearly $\operatorname{tr}(\Psi(a))=$ $r \cdot \operatorname{Tr}_{K / \mathbf{Q}}(a)$, and so by $(2.8), \operatorname{Tr}_{K / \mathbf{Q}}\left(a a^{\rho}\right)>0$ for $a \neq 0$. Then it is an easy exercise to show that $\rho$ is the identity map if $K=F$, and $\rho$ is the Galois involution of $K / F$ if $K \neq F$ (see [S98, p.37, Lemma 2]).
4.2. Put $f(x, y)=E_{X}(q(x), q(y))$. Then $(x, y) \mapsto f(x, y)$ is a $\mathbf{Q}$-valued alternating form on $K_{r}^{1} \times K_{r}^{1}$ such that $f(a x, y)=f\left(x, a^{\rho} y\right)$. For fixed $x, y$ we consider a Q-linear map $a \mapsto f(a x, y)$ of $K$ into $\mathbf{Q}$, and find an element $g(x, y) \in K$ such that $f(a x, y)=\operatorname{Tr}_{K / \mathbf{Q}}(a \cdot g(x, y))$. Then we easily see that $(x, y) \mapsto g(x, y)$ is a skew-hermitian form. Putting $g(x, y)=x \mathcal{T} y^{*}$ with an element $\mathcal{T}$ of $G L_{r}(K)$, where $y^{*}={ }^{t} y^{\rho}$ as we defined in (1.5), we thus obtain

$$
\begin{align*}
\mathcal{T}^{*} & =-\mathcal{T}  \tag{4.5}\\
E_{X}(q(x), q(y)) & =\operatorname{Tr}_{K / \mathbf{Q}}\left(x \mathcal{T} y^{*}\right) \quad \text { for every } \quad x, y \in K_{r}^{1} \tag{4.6}
\end{align*}
$$

With a Z-lattice $L$ in $K_{r}^{1}$ and $\left\{u_{i}\right\}_{i=1}^{s} \subset K_{r}^{1}$ we consider a set of data

$$
\begin{equation*}
\Omega=\left\{K, \Psi, L, \mathcal{T},\left\{u_{i}\right\}_{i=1}^{s}\right\} . \tag{4.7}
\end{equation*}
$$

We call such an $\Omega$ a PEL-type. We note here an easy fact:
(4.8) The direct sum of $\Psi$ and its complex conjugate is equivalent to a $\mathbf{Q}$-rational representation of $K$.
This follows from (2.7). Given such an $\Omega$, we say that $\mathcal{P}=\left(A, \mathcal{C}, \iota ;\left\{t_{i}\right\}_{i=1}^{s}\right)$ is of type $\Omega$ (with respect to $\xi$ and $q$ ) if there is an $\mathbf{R}$-linear isomorphism $q:\left(K_{\mathbf{a}}\right)_{r}^{1}$ $\rightarrow \mathbf{C}^{d}$ with which we have (4.3) and such that $q(a x)=\Psi(a) q(x)$ and $\iota(a) \circ \xi=$ $\xi \circ \Psi(a)$ for every $a \in K, q\left(u_{i}\right)=t_{i}$ for every $i$, and (4.6) holds with some $X \in \mathcal{C}$.
4.3. Let us now classify all structures of type $\Omega$. We first treat Case SP , in which $K=F$. Then ${ }^{t} \mathcal{T}=-\mathcal{T}$. Changing the coordinate system in $F_{r}^{1}$, we may assume that $\mathcal{T}=\eta_{n}$ with $n=r / 2$. This means that $r$ must be even; then $d=n[F: \mathbf{Q}]$. Also any $\mathbf{Q}$-rational representation of $F$ must be equivalent to a multiple of the regular representation of $F$ over $\mathbf{Q}$, and hence (4.8) shows that $\Psi$ must be equivalent to the direct sum of $n$ copies of the regular representation of $F$ over $\mathbf{Q}$. Namely we can decompose $\mathbf{C}^{d}$ into the direct sum $\bigoplus_{v \in \mathbf{a}} V_{v}$ so that each $V_{v}$ is isomorphic to $\mathbf{C}^{n}$ and $\Psi(a)$ acts on $V_{v}$ as a scalar $a_{v}$ for each $a \in F$.

Next take $K$ to be a CM-field, and $\tau=\left\{\tau_{v}\right\}_{v \in \mathbf{a}}$ to be a CM-type as in $\S 3.5$. Let $m_{v}$ resp. $n_{v}$ be the multiplicity of $\tau_{v}$ resp. $\rho \tau_{v}$ in $\Psi$. Then (4.8) shows that
$m_{v}+n_{v}$ must be the same for all $v \in \mathbf{a}$. Then $m_{v}+n_{v}=r$ since $2 d=r[K: \mathbf{Q}]$. This time we can decompose $\mathbf{C}^{d}$ into the direct sum $\bigoplus_{v \in \mathbf{a}} V_{v}$ so that each $V_{v}$ is isomorphic to $\mathbf{C}^{r}$ and $\Psi(a)$ acts on $V_{v}$ as $\operatorname{diag}\left[\bar{a}_{v} 1_{m_{v}}, a_{v} 1_{n_{v}}\right]$ for each $a \in K$.

Putting $l=n$ if $K=F$ and $l=r$ if $K \neq F$, we can thus put

$$
\begin{align*}
\mathbf{C}^{d} & =\left(\mathbf{C}^{l}\right)^{\mathbf{a}}, \quad \Psi(a)=\operatorname{diag}\left[\Psi_{v}(a)\right]_{v \in \mathbf{a}},  \tag{4.9}\\
\Psi_{v}(a) & = \begin{cases}a_{v} 1_{n} & (a \in F=K), \\
\operatorname{diag}\left[\bar{a}_{v} 1_{m_{v}}, a_{v} 1_{n_{v}}\right] & (a \in K \neq F) .\end{cases} \tag{4.10}
\end{align*}
$$

Then (4.3) can be written

with a map $q$ such that

$$
\begin{gather*}
q(a x)=\Psi(a) q(x) \quad\left(a \in K_{\mathbf{a}}, x \in\left(K_{\mathbf{a}}\right)_{r}^{1}\right),  \tag{4.12}\\
E_{X}(q(x), q(y))=\operatorname{Tr}_{K_{\mathbf{a}} / \mathbf{R}}\left(x \mathcal{T} y^{*}\right) \text { for every } x, y \in\left(K_{\mathbf{a}}\right)_{r}^{1},  \tag{4.13}\\
q\left(u_{i}\right)=t_{i} \quad \text { for every } i . \tag{4.14}
\end{gather*}
$$

Conditions (2.1) and (2.3) can be written in the forms

$$
\begin{equation*}
\operatorname{Tr}_{K / \mathbf{Q}}\left(x \mathcal{T} y^{*}\right) \in \mathbf{Z} \quad \text { for every } \quad x, y \in L \tag{4.15}
\end{equation*}
$$

$$
\begin{equation*}
E_{X}(q(x), i \cdot q(y)) \text { is symmetric in }(x, y) \text { and positive definite. } \tag{4.16}
\end{equation*}
$$

Condition (4.15) concerns only $L$, and so we consider $\Omega$ with $L$ satisfying (4.15). We shall later prove the following compatibility condition on $\Psi$ and $\mathcal{T}$ :
(4.17) If a structure of type $\Omega$ exists with $K \neq F$, then the hermitian form $i \mathcal{T}_{v}$ has signature ( $m_{v}, n_{v}$ ) for every $v \in \mathbf{a}$.
Once this is established, we take $Q_{v}$ as in (3.34), consider $\widetilde{G}, G$, and other symbols as in $\S 3.5$ with the present $\mathcal{T}$ and ( $m_{v}, n_{v}$ ). We call this setting Case UB in accordance with what we said in Section 3.

Now, still with a CM-field as $K$, suppose that $\mathcal{T}=\eta_{n}$. Then (4.17) implies that $m_{v}=n_{v}$ for every $v \in \mathbf{a}$. We call this setting Case UT. In this case we put $n=r / 2$; then $m_{v}=n_{v}=n$.
4.4. We are going to determine all $\mathcal{P}$ of type $\Omega$ of (4.7), which amounts to the classification of all possible maps $q$. Let $\left\{e_{k}\right\}_{k=1}^{r}$ be the standard basis of $K_{r}^{1}$. Given any map $q:\left(K_{\mathbf{a}}\right)_{r}^{1} \rightarrow\left(\mathbf{C}^{l}\right)^{\text {a }}$ satisfying (4.12) (but not necessarily injective), we observe that $q$ is determined by the vectors $q\left(e_{k}\right)$, and define a matrix $X_{v}(q) \in \mathbf{C}_{r}^{r}$ for each $v \in \mathbf{a}$ by

$$
\begin{align*}
& X_{v}(q)=\left[\begin{array}{lll}
x_{v}^{1} & \cdots & x_{v}^{2 n} \\
\bar{x}_{v}^{1} & \cdots & \bar{x}_{v}^{2 n}
\end{array}\right], \quad q\left(e_{k}\right)_{v}=x_{v}^{k}  \tag{4.18}\\
& X_{v}(q)=\left[\begin{array}{lll}
x_{v}^{1} & \cdots & x_{v}^{r} \\
\bar{y}_{v}^{1} & \cdots & \bar{y}_{v}^{r}
\end{array}\right], \quad q\left(e_{k}\right)_{v}=\left[\begin{array}{l}
x_{v}^{k} \\
y_{v}^{k}
\end{array}\right] \quad \text { (Case SP), } \tag{4.19}
\end{align*}
$$

Here $x_{v}^{k} \in \mathbf{C}^{n}$ in Case SP and $x_{v}^{k} \in \mathbf{C}^{m_{v}}, y_{v}^{k} \in \mathbf{C}^{n_{v}}$ in Cases UT and UB. We easily see that given any $\omega \in\left(\mathbf{C}_{2 n}^{n}\right)^{\mathbf{a}}$, (resp. $\left.Y \in\left(\mathbf{C}_{r}^{r}\right)^{\mathbf{a}}\right)$ we can find $q$ in Case SP (resp. Cases UT and UB) satisfying (4.12) such that $X_{v}(q)=\left[\begin{array}{c}\omega_{v} \\ \bar{\omega}_{v}\end{array}\right]\left(\right.$ resp. $\left.X_{v}(q)=Y_{v}\right)$ for every $v \in \mathbf{a}$. For the moment we disregard (4.13), (4.15), and (4.16).
4.5. Lemma. (1) For $q$ as above and $\beta \in\left(K_{\mathbf{a}}\right)_{r}^{r}$ put $q^{\beta}(x)=q(x \beta)$ for $x \in$ $\left(K_{\mathbf{a}}\right)_{r}^{1}$. Then $q^{\beta}$ satisfies (4.12) and $X_{v}\left(q^{\beta}\right)=X_{v}(q) \beta_{v}^{*}$.
(2) The map $q$ is injective (and hence surjective) if and only if $\operatorname{det} X_{v}(q) \neq 0$ for every $v \in \mathbf{a}$.

Proof. Assertion (1) can easily be verified. Suppose $q\left(\sum_{i=1}^{r} b_{i} e_{i}\right)=0$ for some $b=\left(b_{i}\right) \in\left(K_{\mathbf{a}}\right)_{r}^{1}$. Let $\beta$ be the element of $\left(K_{\mathbf{a}}\right)_{r}^{r}$ whose rows are all equal to $b$. Then $q^{\beta}\left(e_{k}\right)=q\left(e_{k} \beta\right)=0$ for every $k$, and hence $0=X_{v}\left(q^{\beta}\right)=X_{v}(q) \beta_{v}^{*}$. If $\operatorname{det} X_{v}(q) \neq 0$ for every $v \in \mathbf{a}$, then $\beta_{v}=0$ for every $v \in \mathbf{a}$, that is, $q$ is injective. To prove the converse, we take $q_{0}$ so that $X_{v}\left(q_{0}\right)=\left[\begin{array}{cc}1_{n} & i 1_{n} \\ 1_{n} & -i 1_{n}\end{array}\right]$ in Case SP and $X_{v}\left(q_{0}\right)=1_{r}$ in Cases UT and UB for every $v \in \mathbf{a}$. This is possible by virtue of the observation at the end of $\S 4.4$. Then $q_{0}$ is injective since $X_{v}\left(q_{0}\right)$ is invertible. Given an injective $q, q_{0}^{-1} q$ is a $K_{\mathbf{a}}$-automorphism of $\left(K_{\mathbf{a}}\right)_{r}^{1}$, and so we have $\left(q_{0}^{-1} q\right)(x)=x \beta$ with $\beta \in G L_{r}\left(K_{\mathbf{a}}\right)$. Then $q=q_{0}^{\beta}$, and hence $X_{v}(q)=$ $X_{v}\left(q_{0}\right) \beta_{v}^{*}$, which is invertible. This completes the proof.
4.6. Given an injective $q$, we see that the map $x \mapsto q^{-1}(i \cdot q(x))$ is a $K_{\mathbf{a}^{-}}$ automorphism of $\left(K_{\mathbf{a}}\right)_{r}^{1}$, and hence we have $i \cdot q(x)=q(x C)$ with $C \in G L_{r}\left(K_{\mathbf{a}}\right)$. Then $i \cdot q\left(e_{k}\right)=q^{C}\left(e_{k}\right)$, and so from (4.18), (4.19), and (1) of Lemma 4.5 we obtain

$$
\begin{equation*}
i \cdot I_{v} X_{v}(q)=X_{v}(q) C_{v}^{*} \quad \text { with } \quad I_{v}=\operatorname{diag}\left[1_{m_{v}},-1_{n_{v}}\right] \tag{4.20}
\end{equation*}
$$

where $m_{v}=n_{v}=n$ in Cases SP and UT. We now take (4.13) and (4.16) into account. We have

$$
\begin{aligned}
E_{X}(q(x), i \cdot q(y)) & =E_{X}(q(x), q(y C))=\operatorname{Tr}_{K_{\mathbf{a}} / \mathbf{R}}\left(x \mathcal{T}(y C)^{*}\right) \\
& =[K: F] \sum_{v \in \mathbf{a}} \operatorname{Re}\left(x_{v} \mathcal{T}_{v} C_{v}^{*} y_{v}^{*}\right)
\end{aligned}
$$

This must be symmetric in ( $x, y$ ) and positive definite, which is so if and only if $\mathcal{I}_{v} C_{v}^{*}$ is hermitian and positive definite for every $v \in \mathbf{a}$. Fixing our attention to one $v$, dropping the subscript $v$, and putting simply $X=X_{v}(q)$, from (4.20) we obtain $\mathcal{T} C^{*}=i \mathcal{T} X^{-1} I X$. This is hermitian if and only if $X \mathcal{T}^{-1} X^{*}$ commutes with $I$, that is, if and only if $X \mathcal{T}^{-1} X^{*}$ is of the form $\operatorname{diag}[\beta, \gamma]$ with $\beta \in G L_{m}(\mathbf{C})$ and $\gamma \in G L_{n}(\mathbf{C})$. Then

$$
X\left(\mathcal{T} C^{*}\right)^{-1} X^{*}=-i I X \mathcal{T}^{-1} X^{*}=\operatorname{diag}[-i \beta, i \gamma]
$$

so that both $-i \beta$ and $i \gamma$ are hermitian and positive definite. Thus $i X \mathcal{T}^{-1} X^{*}=$ $\operatorname{diag}[i \beta, i \gamma]$. This proves (4.17). Applying Lemma 3.2 to $X^{*}$ (resp. $Q^{-1} X^{*}$ ), and reinstating the subscript $v$, we thus obtain

$$
X_{v}(q)= \begin{cases}\operatorname{diag}\left[\xi_{v}, \zeta_{v}\right] B\left(z_{v}\right)^{*} & (\text { Cases SP and UT) }  \tag{4.21}\\ \operatorname{diag}\left[\xi_{v}, \zeta_{v}\right] B\left(z_{v}\right) Q_{v}^{*} & (\text { Case UB })\end{cases}
$$

with $\xi_{v} \in G L_{m_{v}}(\mathbf{C}), \zeta_{v} \in G L_{n_{v}}(\mathbf{C})$, and $\left(z_{v}\right) \in \mathcal{H} ; \zeta_{v}=\bar{\xi}_{v}$ in Case SP. Conversely, given $X_{v}(q)$ in this fashion, we see that $q$ is injective, and reversing our
reasoning, we find that $\mathcal{I}_{v} C_{v}^{*}$ is hermitian and positive definite, and hence (4.16) holds.
4.7. Given $z=\left(z_{v}\right) \in \mathcal{H}$, define $p:\left(K_{\mathbf{a}}\right)_{r}^{1} \times \mathcal{H} \rightarrow\left(\mathbf{C}^{l}\right)^{\mathbf{a}}$ and $p_{z}:\left(K_{\mathbf{a}}\right)_{r}^{1} \rightarrow\left(\mathbf{C}^{l}\right)^{\mathbf{a}}$ by specifying $X_{v}\left(p_{z}\right)$ as follows:

$$
\begin{align*}
& X_{v}\left(p_{z}\right)= \begin{cases}B\left(z_{v}\right)^{*} & (\text { Cases SP and UT) }, \\
B\left(z_{v}\right) Q_{v}^{*} & \text { (Case UB), }\end{cases}  \tag{4.22}\\
& p(x, z)=p_{z}(x) \quad\left(x \in\left(K_{\mathbf{a}}\right)_{r}^{1}, z \in \mathcal{H}\right) . \tag{4.23}
\end{align*}
$$

It can easily be seen that $p_{z}(x)$ is holomorphic in $z$. In particular, in Cases SP and UT we have

$$
\left.p_{z}(x)=\left\{\begin{array}{ll}
\left(\left[\begin{array}{ll}
z_{v} & 1_{n}
\end{array}\right] \cdot{ }^{t} x_{v}\right)_{v \in \mathbf{a}} & \left(x \in\left(F_{\mathbf{a}}\right)_{2 n}^{1}, z \in \mathfrak{H}^{\mathbf{a}}\right)  \tag{4.24}\\
\left(\left[\begin{array}{ll}
z_{v} & 1_{n}
\end{array}\right] x_{v}^{*},\left[{ }^{t} z_{v}\right.\right. & 1_{n}
\end{array}\right] \cdot{ }^{t} x_{v}\right)_{v \in \mathbf{a}} \quad\left(x \in\left(K_{\mathbf{a}}\right)_{2 n}^{1}, z \in \mathcal{H}^{\mathbf{a}}\right) .
$$

As observed at the end of $\S 4.6$, from $p_{z}$ we obtain a structure $\mathcal{P}_{z}$ of type $\Omega$ for which (4.11) holds with $p_{z}$ as $q$. To be more precise, by Lemma $4.5(2), p_{z}(L)$ is a lattice in $\left(\mathbf{C}^{l}\right)^{\mathbf{a}}$ so that $\left(\mathbf{C}^{l}\right)^{\mathbf{a}} / p_{z}(L)$ is a complex torus. Define $E_{z}$ by

$$
\begin{equation*}
E_{z}\left(p_{z}(x), p_{z}(y)\right)=\operatorname{Tr}_{K_{\mathbf{a}} / \mathbf{R}}\left(x \mathcal{T} y^{*}\right) \tag{4.25}
\end{equation*}
$$

Since (4.15) and (4.16) are satisfied, $E_{z}$ is a Riemann form, so that $\left(\mathbf{C}^{l}\right)^{\mathbf{a}} / p_{z}(L)$ has a structure of an abelian variety; call it $A_{z}$ and denote by $\mathcal{C}_{z}$ the polarization of $A_{z}$ given by $E_{z}$. For $a \in K$ denote by $\iota_{z}(a)$ the element of $\operatorname{End}_{\mathbf{Q}}\left(A_{z}\right)$ represented by $\Psi(a)$ of (4.9) and (4.10), and by $t_{i}(z)$ the point of $A_{z}$ represented by $p_{z}\left(u_{i}\right)$. Thus we obtain a nonempty family of polarized abelian varieties

$$
\begin{equation*}
\mathcal{F}(\Omega)=\left\{\mathcal{P}_{z} \mid z \in \mathcal{H}\right\}, \quad \mathcal{P}_{z}=\left(A_{z}, \mathcal{C}_{z}, \iota_{z} ;\left\{t_{i}(z)\right\}_{i=1}^{s}\right) \tag{4.26}
\end{equation*}
$$

under the condition, which we hereafter assume, that
(4.27) In Case UB the hermitian form $i \mathcal{T}_{v}$ has signature ( $m_{v}, n_{v}$ ) for every $v \in \mathbf{a}$.
4.8. Theorem. (1) $\mathcal{P}_{z}$ is of type (4.7) for every $z \in \mathcal{H}$.
(2) A structure of type (4.7) is isomorphic to $\mathcal{P}_{z}$ for some $z \in \mathcal{H}$.
(3) $\mathcal{P}_{z}$ and $\mathcal{P}_{w}$ are isomorphic if and only if $w=\gamma z$ for some $\gamma \in \Gamma$, where

$$
\begin{equation*}
\Gamma=\left\{\alpha \in G \mid L \alpha=L \text { and } u_{i} \alpha-u_{i} \in L \text { for every } i\right\} \tag{4.28}
\end{equation*}
$$

Proof. Assertion (1) is obvious. Given $\mathcal{P}$ of type $\Omega$, take $q$ as in (4.11); then we obtain $\xi_{v}, \eta_{v}, \zeta_{v}$, and $z \in \mathcal{H}$ as in (4.21). Define $S \in G L_{l}(\mathbf{C})^{\mathbf{a}}$ by $S=\operatorname{diag}\left[\kappa_{v}\right]_{v \in \mathbf{a}}$ with $\kappa_{v}=\operatorname{diag}\left[\xi_{v}\right]$ in Case SP and $\kappa_{v}=\operatorname{diag}\left[\xi_{v}, \bar{\zeta}_{v}\right]$ in Cases UT and UB. Clearly $S \Psi(a)=\Psi(a) S$ for every $a \in K$. Now (4.21) and (4.22) show that $q\left(e_{k}\right)=S p_{z}\left(e_{k}\right)$ for every $k$, that is, $q=S \circ p_{z}$. We easily see that $S$ gives an isomorphism of $\mathcal{P}_{z}$ onto $\mathcal{P}$, which proves (2). Before proving (3) we make some preliminary observations.

If $\sum_{v \in \mathbf{a}} m_{v} n_{v}=0$ in Case UB, then $\mathcal{H}$ consists of a single point. Thus, under (4.27) there is exactly one isomorphism-class of structures $\mathcal{P}$ of type $\Omega$. In fact, it can be shown that this $\mathcal{P}$ is isogenous to the product of $r$ copies of an abelian variety belonging to a CM-type; see [S98, Theorem 24.15].
4.9. For $\alpha \in \widetilde{G}_{+}$and $z \in \mathcal{H}$ define $M(\alpha, z) \in G L_{l}(\mathbf{C})^{\text {a }}$ by

$$
M(\alpha, z)= \begin{cases}\operatorname{diag}\left[\mu_{v}(\alpha, z)\right]_{v \in \mathbf{a}} & (\text { Case SP) }  \tag{4.29}\\ \operatorname{diag}\left[\lambda_{v}(\alpha, z), \mu_{v}(\alpha, z)\right]_{v \in \mathbf{a}} & \text { (Cases UT and UB) }\end{cases}
$$

Formula (3.14) or (3.38), combined with (4.22) and Lemma 4.5(1), gives

$$
\begin{equation*}
X_{v}\left(p_{z}^{\alpha}\right)=X_{v}\left(p_{z}\right) \alpha_{v}^{*}=\operatorname{diag}\left[{ }^{t} \lambda_{v}(\alpha, z), \mu_{v}(\alpha, z)^{*}\right] X_{v}\left(p_{\alpha z}\right) \tag{4.30}
\end{equation*}
$$

(See $\S 3.3$ for the meaning of the symbols when $m_{v} n_{v}=0$.) This means that $p_{z}^{\alpha}\left(e_{k}\right)=$ ${ }^{t} M(\alpha, z) p_{\alpha z}\left(e_{k}\right)$ for every $k$. Since ${ }^{t} M(\alpha, z)$ commutes with $\Psi(a)$ for every $a \in K$, we have $p_{z}(x \alpha)=p_{z}^{\alpha}(x)={ }^{t} M(\alpha, z) p_{\alpha z}(x)$, that is,

$$
\begin{equation*}
p(x \alpha, z)={ }^{t} M(\alpha, z) p(x, \alpha z) \quad\left(x \in\left(K_{\mathbf{a}}\right)_{r}^{1}, \alpha \in \widetilde{G}_{+}, z \in \mathcal{H}\right) . \tag{4.31}
\end{equation*}
$$

If $\alpha$ belongs to $\Gamma$ of (4.28), then we easily see that ${ }^{t} M(\alpha, z)$ sends $p_{\alpha z}(L)$ onto $p_{z}(L)$, and also sends $E_{z}$ back to $E_{\alpha z}$ by virtue of (4.25). By (2.5) this means that it defines an isomorphism of $\mathcal{P}_{\alpha z}$ onto $\mathcal{P}_{z}$, which proves the "if"-part of Theorem 4.8 (3). Conversely, suppose that there is an isomorphism of $\mathcal{P}_{w}$ to $\mathcal{P}_{z}$; represent it by $S \in G L_{l}(\mathbf{C})^{\mathbf{a}}$. Then $p_{z}^{-1} \circ S \circ p_{w}$ defines an element $\alpha$ of $G L_{r}\left(K_{\mathbf{a}}\right)$, that is, $S p_{w}(x)=p_{z}(x \alpha)$. Since $S$ sends $\left(p_{w}(L), \mathcal{C}_{w}, p_{w}\left(u_{i}\right)\right)$ to $\left(p_{z}(L), \mathcal{C}_{z}, p_{z}\left(u_{i}\right)\right)$, we see that $\alpha \in \Gamma$ by virtue of (2.5). This completes the proof of Theorem 4.8.
4.10. Given a $\mathfrak{g}$-lattice $L$ in $K_{r}^{1}$ and an integral $\mathfrak{g}$-ideal $\mathfrak{c}$, we put

$$
\begin{equation*}
\Gamma(L, \mathfrak{c})=\{\alpha \in \widetilde{G} \mid L \alpha=L \text { and } L(1-\alpha) \subset \mathfrak{c} L\} . \tag{4.32}
\end{equation*}
$$

We call a subgroup $\Gamma$ of $\widetilde{G}$ (resp. $G$ ) a congruence subgroup of $\widetilde{G}$ (resp. $G$ ) if $\Gamma$ contains $\Gamma(L, \mathfrak{c})$ (resp. $\Gamma(L, \mathfrak{c}) \cap G$ ) as a subgroup of finite index for some $L$ and c. Clearly the group of (4.28) is a congruence subgroup of $G$. We note here two easy facts:
(4.33) $\operatorname{det}(\alpha)$ and $\nu(\alpha)$ for $\alpha$ in such a $\Gamma$ are units; $\operatorname{det}(\alpha)$ is a root of unity if $\alpha \in \Gamma \cap G$.

$$
\begin{equation*}
\operatorname{det}(\alpha)=1 \text { if } \alpha \in G \cap \Gamma(L, m \mathfrak{g}) \text { with } 2<m \in \mathbf{Z} \tag{4.34}
\end{equation*}
$$

The first part of (4.33) is obvious. If $\alpha \in \Gamma \cap G$, then $\operatorname{det}(\alpha)$ is a unit in $K$, and $\operatorname{det}(\alpha) \operatorname{det}(\alpha)^{\rho}=1$ by (1.12), and so $|\operatorname{det}(\alpha)|_{v}=1$ for every $v \in \mathbf{a}$ by (1.22). Therefore $\operatorname{det}(\alpha)$ is a root of unity. If $\alpha \in \Gamma(L, m \mathfrak{g})$, then $\operatorname{det}(\alpha)-1$ is divisible by $m$, and hence we obtain (4.34).

It is well-known that for a congruence subgroup $\Gamma$ of $G$ the quotient space $\Gamma \backslash \mathcal{H}$ has a compactification, called the Satake compactification, which has a structure of complex analytic space, and as such, is isomorphic to a normal projective variety $V^{*}$; moreover, $\Gamma \backslash \mathcal{H}$ is mapped onto a Zariski open subset $V$ of $V^{*}$. Let $\varphi$ denote the $\Gamma$-invariant map $\mathcal{H} \rightarrow V$ that gives this isomorphism. We then call $(V, \varphi)$ a model of $\Gamma \backslash \mathcal{H}$.
4.11. We now introduce the notion of $C M$-points on $\mathcal{H}$. We take a CM-algebra $Y=K_{1} \oplus \cdots \oplus K_{t}$ with CM-fields $K_{i}$ as in $\S 2.9$. We assume that $K \subset K_{i}$ for every $i$ and $r=[Y: K]$. We denote by $F_{i}$ the maximal real subfield of $K_{i}$, and by $\rho$ the automorphism of $Y$ that coincides with the Galois involution of $K_{i} / F_{i}$ for every $i$. Let us now consider a $K$-linear ring-injection $h: Y \rightarrow K_{r}^{r}$ satisfying

$$
\begin{equation*}
h\left(a^{\rho}\right)=\mathcal{T} h(a)^{*} \mathcal{T}^{-1} \quad(a \in Y) \tag{4.35}
\end{equation*}
$$

Put $Y^{u}=\left\{a \in Y \mid a a^{\rho}=1\right\}$. Then clearly $h\left(Y^{u}\right) \subset G$. Since $Y^{u}$ is contained in a compact subgroup of $\left(Y \otimes_{\mathbf{Q}} \mathbf{R}\right)^{\times}$, we easily see that the projection of $h\left(Y^{u}\right)$ to $G_{\mathbf{a}}$ is contained in a compact subgroup of $G_{\mathbf{a}}$, and hence $h\left(Y^{u}\right)$ has a common fixed point in $\mathcal{H}$. Moreover, $h\left(Y^{u}\right)$ has only one common fixed point as will be shown in

Lemma 4.12 below. We call a point on $\mathcal{H}$ which is obtained as such a fixed point a $C M$-point on $\mathcal{H}$ with rèspect to $G$. If $w$ is a CM-point and $\beta \in \widetilde{G}_{+}$, then $\beta(w)$ is also a CM-point, since $h^{\prime}$ defined by $h^{\prime}(a)=\beta h(a) \beta^{-1}$ satisfies (4.35) and $\beta(w)$ is fixed by $h^{\prime}\left(Y^{u}\right)$.

Let us now show that in Case SP an injection $h$ of type (4.35) always exists for any given $Y$. Take an element $\zeta$ of $Y^{\times}$such that $\zeta^{\rho}=-\zeta$. Then $(x, y) \mapsto$ $\operatorname{Tr}_{Y / F}\left(\zeta x y^{\rho}\right)$ is an $F$-valued nondegenerate alternating form on $Y \times Y$. Therefore we can find an $F$-linear bijection $r: Y \rightarrow F_{2 n}^{1}$ such that $\operatorname{Tr}_{Y / F}\left(\zeta x y^{\rho}\right)=r(x) \eta_{n} \cdot{ }^{t} r(y)$. Since multiplication by $a \in Y$ is an $F$-linear endomorphism of $Y$, we can define an $F$-linear map $h: Y \rightarrow F_{2 n}^{2 n}$ by the relation $r(a x)=r(x) h(a)$ for every $a, x \in Y$. Then we can easily verify (4.35).

In the unitary case the matter is not so simple. However, we can find at least one ( $Y, h$ ) as follows. We first consider Case UB. Changing the coordinate system of $K_{r}^{1}$, we may assume that $\mathcal{T}$ is diagonal. Take $Y$ to be the direct sum of $r$ copies of $K$ and define $h$ by $h(a)=\operatorname{diag}\left[a_{1}, \ldots, a_{r}\right]$ for $a=\left(a_{i}\right)_{i=1}^{r} \in Y$ with $a_{i} \in K$. Clearly (4.35) is satisfied. In Case UT we can find $\sigma \in G L_{2 n}(K)$ such that $\sigma \eta_{n} \sigma^{*}$ is diagonal, and so the argument in Case UB produces a map in Case UT. Thus $\mathcal{H}$ has at least one CM-point with respect to $G$. Also we can always find CM-points on $\mathcal{H}$ associated with infinitely many different CM-fields $Y$, as shown in [S64, Proposition 4.10] and [S66b, pp.379-381].

Now let $w$ denote the CM-point obtained as the common fixed point of $h\left(Y^{u}\right)$ as above in all three cases. Putting $X_{v}=X_{v}\left(p_{w}\right)$, from (4.30) and (4.31) we obtain

$$
\begin{array}{lll}
h(\alpha)_{v} X_{v}^{*}=X_{v}^{*} \operatorname{diag}\left[\overline{\lambda_{v}(h(\alpha), w)}, \mu_{v}(h(\alpha), w)\right] & (v \in \mathbf{a}), \\
p_{w}(x h(\alpha))=^{t} M(h(\alpha), w) p_{w}(x) & \left(x \in\left(K_{\mathbf{a}}\right)_{r}^{1}, \alpha \in Y^{u}\right) \tag{4.36b}
\end{array}
$$

Since $Y^{u}$ spans $Y$ over $\mathbf{Q}$ as will be shown in Lemma 4.12 below, we can extend these equalities $\mathbf{Q}$-linearly to $Y$ and define $\psi_{v}: Y \rightarrow \mathbf{C}_{m_{v}}^{m_{v}}, \varphi_{v}: Y \rightarrow \mathbf{C}_{n_{v}}^{n_{v}}$, and $\Phi: Y \rightarrow \operatorname{End}\left(\left(\mathbf{C}^{n}\right)^{\mathbf{a}}\right)$ so that

$$
\begin{align*}
& \psi_{v}(\alpha)=\lambda_{v}(h(\alpha), w) \text { and } \varphi_{v}(\alpha)=\mu_{v}(h(\alpha), w) \quad \text { for } \quad \alpha \in Y^{u},  \tag{4.37}\\
& h(a)_{v} X_{v}^{*}=X_{v}^{*} \operatorname{diag}\left[\overline{\psi_{v}(a)}, \varphi_{v}(a)\right] \quad(a \in Y, v \in \mathbf{a}),  \tag{4.38}\\
& p_{w}(x h(a))={ }^{t} \Phi(a) p_{w}(x) \quad\left(x \in\left(K_{\mathbf{a}}\right)_{r}^{1}, a \in Y\right),  \tag{4.39}\\
& \Phi(a)=\operatorname{diag}\left[\Phi_{v}(a)\right]_{v \in \mathbf{a}}, \quad \Phi_{v}(a)= \begin{cases}\varphi_{v}(a) & (K=F), \\
\operatorname{diag}\left[\psi_{v}(a), \varphi_{v}(a)\right] & (K \neq F) .\end{cases} \tag{4.40}
\end{align*}
$$

(In Case SP we have $\psi_{v}=\varphi_{v}$ and $m_{v}=n_{v}=n$.) Notice that the last equality follows from (4.29). From (4.38) and (4.40) we see that $\Phi_{v}(a)=\Psi_{v}(a)$ for every $a \in K$. Therefore the restriction of $\Phi$ to $K$ is $\Psi$. Clearly ${ }^{t} \Phi(\alpha)$ for each $\alpha \in Y$ defines an element of $\operatorname{End}_{\mathbf{Q}}\left(A_{w}\right)$; denote it by $\iota^{\prime}(\alpha)$. Then $\iota^{\prime}$ is a ring-injection of $Y$ into $\operatorname{End}_{\mathbf{Q}}\left(A_{w}\right)$ that coincides with $\iota_{w}$ on $K$. Thus we find that $\mathcal{P}_{w}$, together with $\iota^{\prime}$, defines a structure considered in $\S 2.9$ with $Y$ and ${ }^{t} \Phi$ as $W$ and $\Phi$ there, and obtain a CM-type ( $K_{i}, \Phi_{i}$ ) for each $i$ such that $\Phi$ is equivalent to the direct sum of $\Phi_{1}, \ldots, \Phi_{t}$ in the sense of (2.12). Moreover, from (4.25), (4.35), and (4.39) we see that $E_{w}\left({ }^{t} \Phi(\alpha) u, v\right)=E_{w}\left(u,{ }^{t} \Phi\left(\alpha^{\rho}\right) v\right)$. Thus the automorphism $\rho$ of $Y$ corresponds to the involution of $\operatorname{End}_{\mathbf{Q}}\left(A_{w}\right)$ determined by $\mathcal{C}_{w}$ in the sense of $\S 2.5$.
4.12. Lemma. Let $(Y, h)$ be as in $\S 4.11$. Then $Y$ is spanned by $Y^{u}$ over $\mathbf{Q}$, and there exists an element $\beta$ of $Y^{u}$ such that $Y=\mathbf{Q}[\beta]$. Moreover, $h(\beta)$ for any such $\beta$ has only one fixed point in $\mathcal{H}$.

Proof. Clearly $Y \otimes_{\mathbf{Q}} \mathbf{R}$ as an $\mathbf{R}$-algebra is isomorphic to $\mathbf{C}^{d}$. For $y \in Y \otimes_{\mathbf{Q}} \mathbf{R}$ denote by $y_{i}$ the $i$-th coordinate of $y$ viewed as an element of $\mathbf{C}^{d}$. Then $\left(y_{i}\right)^{\rho}=$ $\left(y^{\rho}\right)_{i}$. Since $Y$ is dense in $\mathbf{C}^{d}$, we can find an element $x$ of $Y^{\times}$such that $x_{i} / x_{j} \notin \mathbf{R}$ for $i \neq j$, and $x_{i}^{\rho} / x_{j} \notin \mathbf{R}$ for every $(i, j)$. Put $\beta=x^{\rho} / x$. Then $\beta \in Y^{u}$, and $\beta_{1}, \ldots, \beta_{d}, \beta_{1}^{\rho}, \ldots, \beta_{d}^{\rho}$ are all different. Thus $Y=\mathbf{Q}[\beta]$, and hence $Y$ is spanned by $Y^{u}$ over $\mathbf{Q}$, since the powers of $\beta$ belong to $Y^{u}$. To prove the uniqueness of the fixed point of $h(\beta)$, we first consider Case SP. Employing the symbols of $\S 3.6$, observe that

$$
\begin{equation*}
\left\{\alpha \in G^{\prime} \mid \alpha(0)=0\right\}=\{\operatorname{diag}[u, \bar{u}] \mid u \in U(n)\}, \tag{4.41}
\end{equation*}
$$

where $U(n)=U(n, 0)$ in the sense of (3.3). Let $w$ be a fixed point of $h\left(Y^{u}\right)$ as in $\S 4.11$ in Case SP. Then for each $v \in \mathbf{a}$ we can find $\xi_{v} \in \mathfrak{E}^{\prime}$ such that $\xi_{v}(0)=w_{v}$ and

$$
\begin{equation*}
\xi_{v}^{-1} h(\alpha)_{v} \xi_{v}=\operatorname{diag}\left[\sigma_{v}(\alpha), \bar{\sigma}_{v}(\alpha)\right] \quad\left(\alpha \in Y^{u}\right) \tag{4.42}
\end{equation*}
$$

with a map $\sigma_{v}: Y^{u} \rightarrow U(n)$. Suppose $w^{\prime}$ is another fixed point of $h\left(Y^{u}\right)$ on $\mathcal{H}$; put $z_{v}^{\prime}=\xi_{v}^{-1}\left(w_{v}^{\prime}\right)$. Then $z_{v}^{\prime}=\sigma_{v}(\alpha) z_{v}^{\prime} \cdot{ }^{t} \sigma_{v}(\alpha)$ for every $\alpha \in Y^{u}$. Let $c_{v 1}, \ldots, c_{v n}$ be the characteristic roots of $\sigma_{v}(\alpha)$. Diagonalizing $\sigma_{v}(\alpha)$, we easily see that $z^{\prime}$ must be 0 if $c_{v j} c_{v k} \neq 1$ for every $(j, k)$. Now take $\alpha$ to be the above $\beta$. From (4.42) we see that $c_{v 1}, \ldots, c_{v n}, \bar{c}_{v 1}, \ldots, \bar{c}_{v n}$ are the characteristic roots of $h(\beta)_{v}$, so that $c_{v j} \neq \bar{c}_{v k}$ for every ( $j, k$ ). Thus $c_{v j} c_{v k} \neq 1$ for every ( $j, k$ ) and every $v \in \mathbf{a}$. Then $w_{v}^{\prime}=w_{v}$, which proves the desired fact. Case UT can be handled by the same technique. In Case UB we take an element $\beta$ of $G_{\mathbf{a}}$ so that $\beta(0)=w$. Then taking $\mathfrak{B}\left(m_{v}, n_{v}\right)$ itself in place of $\mathfrak{B}_{n, n}$, we can prove the uniqueness of the fixed point in the same manner.
4.13. Lemma. If $w$ is a CM-point on $\mathcal{H}$, then the entries of $w_{v}, X_{v}\left(p_{w}\right)$ and $\psi_{v}(a), \varphi_{v}(a)$ of (4.38) are algebraic for every $v \in \mathbf{a}$ and every $a \in Y$.

Proof. The point $w$ can be obtained as the unique fixed point of $h(\alpha)$ with some $\alpha \in Y^{u}$ as above. Since we took $Q_{v}$ of (3.34) to be algebraic, the action of $h(\alpha)$ on $\mathcal{H}$ is $\overline{\mathbf{Q}}$-rational. (Notice that $\mathcal{H}$ has an obvious $\overline{\mathbf{Q}}$-rational structure.) Therefore the fixed point $w_{v}$ of $Q_{v}^{-1} h(\alpha)_{v} Q_{v}$ has algebraic entries. Then from (4.22) and (4.38) we immediately see the algebraicity of the other quantities in question.

If $\widetilde{G}=G p(1, F)=G L_{2}(F)$, then $\mathcal{H}=\mathfrak{H}_{1}^{\text {a }}$. This is the so-called Hilbert modular case. In this case the CM-points on $\mathcal{H}$ can be described in a clear-cut way as follows:
4.14. Proposition. Let $\widetilde{G}=G p(1, F)=G L_{2}(F)$. Let $\left(K,\left\{\tau_{v}\right\}_{v \in \mathbf{a}}\right)$ be a $C M-$ type with $K$ containing $F$ as the maximal real subfield and an injection $\tau_{v}: K \rightarrow \mathbf{C}$ which extends $v: F \rightarrow \mathbf{R}$. Take an element $w_{0}$ of $K$ so that $\operatorname{Im}\left(\tau_{v}\left(w_{0}\right)\right)>0$ for every $v \in \mathbf{a}$, and put $w=\left(w_{v}\right)_{v \in \mathbf{a}}$ with $w_{v}=\tau_{v}\left(w_{0}\right)$. Then $w$ is a CM-point on $\mathfrak{H}_{1}^{\mathrm{a}}$ with respect to $G L_{2}(F), \tau_{v}$ coincides with $\varphi_{v}$ of (4.38), and $\left(A_{w}, \iota^{\prime}\right)$ is of type $\left(K,\left\{\tau_{v}\right\}_{v \in \mathbf{a}}\right)$. Conversely every CM-point on $\mathfrak{H}_{1}^{\mathbf{a}}$ can be obtained in such a manner.

Proof. Given such a CM-type and $w_{0}$, we can define an $F$-linear ringinjection $h: K \rightarrow F_{2}^{2}$ by $h(\alpha)\left[\begin{array}{c}w_{0} \\ 1\end{array}\right]=\left[\begin{array}{c}w_{0} \\ 1\end{array}\right] \alpha$ for $\alpha \in K$. (Notice that (4.35) with $\mathcal{T}=\eta_{1}$ is satisfied by every $F$-linear ring-injection $h$ of $K$ into $F_{2}^{2}$.) Taking the image under $\tau_{v}$, we obtain (4.38) for every $\alpha \in K$ with $\tau_{v}$ in place of $\varphi_{v}$, which shows that $w$ is a fixed point of $h\left(K^{\times}\right)$and $\left(A_{w}, \iota^{\prime}\right)$ is of type $\left(K,\left\{\tau_{v}\right\}_{v \in \mathbf{a}}\right)$. Thus $w$ is a CM-point. Conversely, consider a CM-point $w$ obtained from an Flinear ring-injection $h: K \rightarrow F_{2}^{2}$. Then we obtain maps $\varphi_{v}: K \rightarrow \mathbf{C}$ satisfying (4.38). Take $\alpha \in K$ so that $h(\alpha)=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right]$ with $c \neq 0$. Then (4.38) shows that $c_{v} w_{v}+d_{v}=\varphi_{v}(\alpha)$. Put $w_{0}=c^{-1}(\alpha-d)$. Then $w_{v}=\varphi_{v}\left(w_{0}\right)$. This proves the converse part.
4.15. Remark. It should be noted that our formulation is essentially the same as in [S98] and some earlier papers of the author, [S63], [S65], and [S79], for example, but some symbols are defined differently. In particular, in those articles we set $\Psi_{v}(a)=\operatorname{diag}\left[a_{v} 1_{m_{v}}, \bar{a}_{v} 1_{n_{v}}\right]$ for $a \in K \neq F$, and formulas (3.34), (3.35), (3.37), and (3.38) were given in accordance with this change; see [S98, §§23 and 24]. Also, the families associated with more general PEL-types are treated in [S63].

## 5. Definition of automorphic forms

5.1. Our setting is the same as in $\S 3.5$; thus $K=F$ and $\mathfrak{r}=\mathfrak{g}$ in Case SP and $K \neq F$ in Cases UT and UB. We now define a symbol $\mathbf{b}$ as follows. In Case SP we put $\mathbf{b}=\mathbf{a}$. In Cases UT and UB we denote by $\mathbf{b}$ the set of all isomorphic embeddings of $K$ into $\mathbf{C}$, and view $\mathbf{a}$ as a subset of $\mathbf{b}$. In all cases, for each $v \in \mathbf{b}$ and $\sigma \in \operatorname{Gal}(\overline{\mathbf{Q}} / \mathbf{Q})$, or more generally for $\sigma \in \operatorname{Aut}(\mathbf{C})$, we denote by $v \sigma$ the element of $\mathbf{b}$ that is the composed map of $v$ and $\sigma$. Then $\mathbf{b}=\mathbf{a} \rho \cup \mathbf{a}$ in Cases UT and UB, where $\rho$ is complex conjugation.

In (3.36) and (3.37) we defined factors of automorphy $\lambda_{v}$ and $\mu_{v}$ for $v \in \mathbf{a}$. We now put, for $\alpha \in \widetilde{G}_{\mathbf{A}+}$ and $z=\left(z_{v}\right)_{v \in \mathbf{a}} \in \mathcal{H}$,

$$
\begin{align*}
& M_{\alpha}(z)=M(\alpha, z)=\left(\mu_{v}(\alpha, z)\right)_{v \in \mathbf{b}}  \tag{5.1}\\
& \mu_{v \rho}(\alpha, z)=\lambda_{v}(\alpha, z), \quad n_{v \rho}=m_{v} \quad(v \in \mathbf{a}, K \neq F), \\
& j_{\alpha}(z)=j(\alpha, z)=\left(j_{v}(\alpha, z)\right)_{v \in \mathbf{b}}, \quad j_{v}(\alpha, z)=\operatorname{det}\left[\mu_{v}(\alpha, z)\right] \quad(v \in \mathbf{b}) . \tag{5.3}
\end{align*}
$$

The symbol $M_{\alpha}(z)$ is the same as $M(\alpha, z)$ of (4.29); in Case SP we shall write also $\mu_{\alpha}(z)$ for $M_{\alpha}(z)$. In Cases UT and UB, we hereafter use $\left(\mu_{v}\right)_{v \in \mathbf{b}}$ and $\left(n_{v}\right)_{v \in \mathbf{b}}$ instead of $\left(\lambda_{v}, \mu_{v}\right)_{v \in \mathbf{a}}$ and $\left(m_{v}, n_{v}\right)_{v \in \mathbf{a}}$. According to the convention of $\S 3.3$ in Case UB , for $v \in \mathbf{a}$, the pair $\left(j_{v \rho}, j_{v}\right)$ is either $(\operatorname{det}(\bar{\alpha}), 1)$ or $(1, \operatorname{det}(\alpha))$ according as $n_{v}=0$ or $n_{v \rho}=0$.

To simplify our notation, for $x, y \in \mathbf{C}^{\mathbf{b}}$ and $\kappa \in \mathbf{C}$ we put

$$
\begin{gather*}
x^{y}=\prod_{v \in \mathbf{b}} x_{v}^{y_{v}}  \tag{5.4a}\\
x^{\kappa \mathbf{a}}=\prod_{v \in \mathbf{a}} x_{v}^{\kappa}, \quad x^{\kappa \mathbf{b}}=\prod_{v \in \mathbf{b}} x_{v}^{\kappa} \tag{5.4~b}
\end{gather*}
$$

The factors $x_{v}^{y_{v}}$ and $x_{v}^{\kappa}$ must be understood according to the context. If $0<x_{v} \in$ $\mathbf{R}$, we always put $x_{v}^{y_{v}}=\exp \left(y_{v} \log x_{v}\right)$ with real $\log x_{v}$. If we identify ка (resp. $\kappa \mathbf{b}$ ) with the element of $\mathbf{Z}^{\mathbf{a}}$ (resp. $\mathbf{Z}^{\mathbf{b}}$ ) whose components are all equal to $\kappa$, then (5.4b) is a special case of (5.4a). We shall later speak of $x^{\kappa \mathbf{a}+\lambda}$ with $\lambda \in \mathbf{C}^{\mathbf{a}}$.

To define automorphic forms on $\mathcal{H}$, we naturally assume that $\operatorname{dim}(\mathcal{H})>0$. Then we take a rational representation

$$
\begin{equation*}
\omega: \prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C}) \rightarrow G L(X) \tag{5.5}
\end{equation*}
$$

with a finite-dimensional complex vector space $X$. We understand that $n_{v}=n$ in Cases SP and UT, and $G L_{n_{v}}(\mathbf{C})=1$ if $n_{v}=0$ in Case UB (see §3.3). Notice that $\prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C})=\prod_{v \in \mathbf{a}}\left[G L_{m_{v}}(\mathbf{C}) \times G L_{n_{v}}(\mathbf{C})\right]$ in Cases UT and UB. In all cases $M_{\alpha}(z)$ is an element of $\prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C})$. We shall often write an element of $\alpha \in \prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C})$ in the form $\alpha=(a, b)$ according to the following convention: $a=b \in G L_{n}(\mathbf{C})^{\mathbf{a}}$ in Case SP; $a=\left(\alpha_{v \rho}\right)_{v \in \mathbf{a}} \in \prod_{v \in \mathbf{a}} G L_{m_{v}}(\mathbf{C})$ and $b=\left(\alpha_{v}\right)_{v \in \mathbf{a}} \in$ $\prod_{v \in \mathbf{a}} G L_{n_{v}}(\mathbf{C})$ in Cases UT and UB.

Given a map $f: \mathcal{H} \rightarrow X$ and $\alpha \in \widetilde{G}_{\mathbf{A +}}$, we define $f \|_{\omega} \alpha: \mathcal{H} \rightarrow X$ and $\left.f\right|_{\omega} \alpha: \mathcal{H} \rightarrow X$ by

$$
\begin{align*}
& \left(f \|_{\omega} \alpha\right)(z)=\omega\left(M_{\alpha}(z)\right)^{-1} f(\alpha z) \quad(z \in \mathcal{H})  \tag{5.6a}\\
& \left.\quad f\right|_{\omega} \alpha=f \|_{\omega}\left(\nu(\alpha)_{\mathbf{a}}^{-1 / 2} \alpha\right) \tag{5.6b}
\end{align*}
$$

where $\nu(\alpha)_{\mathbf{a}}=\left(\nu(\alpha)_{v}\right)_{v \in \mathbf{a}}$. We easily see that $\left(f \|_{\omega} \alpha\right)\left\|_{\omega} \beta=f\right\|_{\omega}(\alpha \beta)$ and $\left.\left(\left.f\right|_{\omega} \alpha\right)\right|_{\omega} \beta$ $=\left.f\right|_{\omega}(\alpha \beta)$.

Given $k=\left(k_{v}\right)_{v \in \mathbf{b}} \in \mathbf{Z}^{\mathbf{b}}$, we can define a representation $\omega: \prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C}) \rightarrow$ $G L(\mathbf{C})=\mathbf{C}^{\times}$by $\omega(x)=\operatorname{det}(x)^{k}$ using the notation of (5.4a). Then we write $f \|_{k} \alpha$ for $f \|_{\omega} \alpha$; thus

$$
\begin{equation*}
\left(f \|_{k} \alpha\right)(z)=j_{\alpha}(z)^{-k} f(\alpha z) \quad(z \in \mathcal{H}) \tag{5.7}
\end{equation*}
$$

We shall often write $f \| \alpha$ for $f \|_{\omega} \alpha$ or $f \|_{k} \alpha$ when $\omega$ or $k$ is clear from the context. We use $f \mid \alpha$ in a similar sense.
5.2. Given a congruence subgroup $\Gamma$ of $\widetilde{G}$ or $G$ contained in $\widetilde{G}_{+}$and $\omega$ as in (5.5), we denote by $\mathcal{M}_{\omega}(\Gamma)$ the set of all functions $f: \mathcal{H} \rightarrow X$ satisfying the following conditions:
(5.8) $f$ is holomorphic;
(5.9) $\left.f\right|_{\omega} \gamma=f$ for every $\gamma \in \Gamma$;
(5.10) $f$ is holomorphic at every cusp.

If $\omega(x)=\operatorname{det}(x)^{k}$ as above, we write $\mathcal{M}_{k}(\Gamma)$ for $\mathcal{M}_{\omega}(\Gamma)$. In particular $\mathcal{M}_{\kappa \mathbf{a}}(\Gamma)$ and $\mathcal{M}_{\kappa \mathbf{b}}(\Gamma)$ are meaningful for $\kappa \in \mathbf{Z}$. Condition (5.10) is necessary or meaningful only in the following exceptional cases: (Cases SP and UT) $F=\mathbf{Q}$ and $n=1$; (Case UB) $F=\mathbf{Q}, r=2$, and $x^{*} \mathcal{T} x=0$ for some $x \in K^{2}, \neq 0$. In these cases $\Gamma$ is commensurable with a "conjugate" of $S L_{2}(\mathbf{Z})$. The precise meaning of (5.10) will be explained in $\S 5.6$ below. In our later treatment, we will have to prove that certain functions are elements of $\mathcal{M}_{\omega}(\Gamma)$. In order to do so, we have to verify (5.10) in those exceptional cases. However, since the verification is always easy and we are mainly interested in the higher-dimensional case, we will not give the proof of the fact on each occasion, leaving the task to the reader.

If $\Gamma \subset G$, we can take $f \|_{\omega} \gamma$ instead of $\left.f\right|_{\omega} \gamma$ in (5.9), but there is a natural example of $\Gamma$ not contained in $G$, for which (5.9) is the right condition. For example, we can take $\Gamma=G L_{2}(\mathfrak{g}) \cap \widetilde{G}_{+}$in the Hilbert modular case. In the present book, however, we consider almost exclusively $\Gamma$ contained in $G$ and also $f \|_{\omega} \gamma$ instead of $\left.f\right|_{\omega} \gamma$. The only exception is Theorem 10.4, in which a group $\Gamma$ not contained in $G$ and the symbol $\left.f\right|_{\omega} \gamma$ appear.

An element of $\mathcal{M}_{\omega}(\Gamma)$ is called a (holomorphic) automorphic form of weight $\omega$ (or, of weight $k$, if $\left.\omega(x)=\operatorname{det}(x)^{k}\right)$ with respect to $\Gamma$. An automorphic form is also called a modular form usually in Case SP.
5.3. Since it is often convenient not to specify $\Gamma$, we denote by $\mathcal{M}_{\omega}$ (resp. $\mathcal{M}_{k}$ ) the union of $\mathcal{M}_{\omega}(\Gamma)$ (resp. $\mathcal{M}_{k}(\Gamma)$ ) for all congruence subgroups $\Gamma$ of $G$, and put

$$
\begin{align*}
\mathcal{A}_{\omega} & =\bigcup_{e}\left\{g^{-1} f \mid f \in \mathcal{M}_{\tau_{e}}, 0 \neq g \in \mathcal{M}_{e}\right\}  \tag{5.11}\\
\mathcal{A}_{\omega}(\Gamma) & =\left\{h \in \mathcal{A}_{\omega} \mid h \|_{\omega} \gamma=h \text { for every } \gamma \in \Gamma\right\} \quad(\Gamma \subset G), \tag{5.12}
\end{align*}
$$

where $e$ runs over $\mathbf{Z}^{\mathbf{b}}$, and $\tau_{e}$ denotes the representation defined by $\tau_{e}(x)=$ $\operatorname{det}(x)^{e} \omega(x)$. If $\omega(x)=\operatorname{det}(x)^{k}$, we denote these by $\mathcal{A}_{k}$ and $\mathcal{A}_{k}(\Gamma)$.
5.4. If $\alpha \in G_{\mathbf{A}}$ in Cases UT and UB, from (3.23) we obtain $j_{v \rho}(\alpha, z)=$ $\operatorname{det}(\alpha)_{v}^{-1} j_{v}(\alpha, z)$ for every $v \in \mathbf{a}$. Therefore, for $k \in \mathbf{Z}^{\mathbf{b}}$ we have

$$
\begin{equation*}
j_{\alpha}(z)^{k}=\prod_{v \in \mathbf{a}} \operatorname{det}(\alpha)_{v}^{-k_{v \rho}} j_{v}(\alpha, z)^{k_{v \rho}+k_{v}} \quad \text { if } \quad \alpha \in G_{\mathbf{A}} \tag{5.13}
\end{equation*}
$$

This means that $j_{\alpha}(z)^{k}=\operatorname{det}(\alpha)^{p} j_{\alpha}(z)^{q}$ with $p, q \in \mathbf{Z}^{\mathbf{a}}$, and so $f \|_{k} \alpha$ and $\mathcal{M}_{k}(\Gamma)$ can be defined in terms of $p$ and $q$ instead of $k$. (That is what we did in [S97, $\S 10.4]$.) Also, by (4.34), $\operatorname{det}(\alpha)=1$ if $\alpha$ belongs to a sufficiently small congruence subgroup. Therefore $\mathcal{M}_{k}=\mathcal{M}_{l}$ if $k_{v \rho}+k_{v}=l_{v \rho}+l_{v}$ for every $v \in \mathbf{a}$ such that the $v$-factor of $\mathcal{H}$ is nontrivial; in particular $\mathcal{M}_{\kappa \mathbf{b}}=\mathcal{M}_{2 \kappa \mathbf{a}}$. In that sense our definition of $\mathcal{M}_{k}$ with $k$ in $\mathbf{Z}^{\mathrm{b}}$ may look awkward, but we shall later see that this is a natural definition from an arithmetical viewpoint. However, since $\mathcal{M}_{k}=\mathcal{M}_{q}$ with some $q \in \mathbf{Z}^{\mathbf{a}}$, we can restrict $e$ to $\mathbf{Z}^{\mathbf{a}}$ in (5.11).
5.5. An element of $\mathcal{A}_{0}(\Gamma)$ is a $\Gamma$-invariant meromorphic function on $\mathcal{H}$, and it is known that conversely every $\Gamma$-invariant meromorphic function on $\mathcal{H}$ belongs to $\mathcal{A}_{0}(\Gamma)$ if we exclude the exceptional cases mentioned in §5.2. An element of $\mathcal{A}_{0}(\Gamma)$ is called an automorphic function with respect to $\Gamma$ or a $\Gamma$-automorphic function on $\mathcal{H}$.

Now, let $(V, \varphi)$ be a model of $\Gamma \backslash \mathcal{H}$ in the sense of $\S 4.10$, and let $\mathbf{C}(V)$ be the field of all functions on $V$ in the sense of algebraic geometry, as defined in §2.4. Then $\mathcal{A}_{0}(\Gamma)$ consists of the functions $g \circ \varphi$ for all $g \in \mathbf{C}(V)$. In this sense $\mathcal{A}_{0}(\Gamma)$ can be identified with $\mathbf{C}(V)$ if we identify $\Gamma \backslash \mathcal{H}$ with $V$.
5.6. Hereafter until the end of Section 8 we confine ourselves to Cases SP and UT; we shall return to Case UB in Sections 9 and 11. In Case UT we identify $\left(\mathbf{C}^{n}\right)^{\mathbf{b}}$ with $\left(\mathbf{C}^{n} \times \mathbf{C}^{n}\right)^{\mathbf{a}}$ through the map $\left(x_{v}\right)_{v \in \mathbf{b}} \mapsto\left(x_{v \rho}, x_{v}\right)_{v \in \mathbf{a}}$, where $x_{v} \in \mathbf{C}^{n}$. Notice that (5.5) becomes $\omega: G L_{n}(\mathbf{C})^{\mathbf{b}} \rightarrow G L(X)$ in Cases SP and UT. We now put

$$
\begin{align*}
& \mathbf{e}(c)=\exp (2 \pi i c) \quad(c \in \mathbf{C})  \tag{5.14}\\
& \mathbf{e}_{\mathbf{a}}(x)=\exp \left(2 \pi i \sum_{v \in \mathbf{a}} x_{v}\right) \quad\left(x \in \mathbf{C}^{\mathbf{a}}\right),  \tag{5.15}\\
& \mathbf{e}_{\mathbf{a}}^{n}(X)=\mathbf{e}_{\mathbf{a}}(\operatorname{tr}(X)) \quad\left(X \in\left(\mathbf{C}_{n}^{n}\right)^{\mathbf{a}}\right),  \tag{5.16}\\
& S=S^{n}=\left\{\sigma \in K_{n}^{n} \mid \sigma^{*}=\sigma\right\} \tag{5.17}
\end{align*}
$$

Observe that $\left[\begin{array}{cc}1 & \sigma \\ 0 & 1\end{array}\right] \in G$ for every $\sigma \in S$ and $\operatorname{diag}[a, \widehat{a}] \in G$ for every $a \in$ $G L_{n}(K)$, and for a function $f: \mathcal{H} \rightarrow X$ we have

$$
\begin{gather*}
\left(f \|_{\omega}\left[\begin{array}{ll}
1 & \sigma \\
0 & 1
\end{array}\right]\right)(z)=f(z+\sigma)  \tag{5.18}\\
\left(f \|_{\omega} \operatorname{diag}[a, \widehat{a}]\right)=\omega\left({ }^{t} a, a^{*}\right) f\left(a z a^{*}\right) \tag{5.19}
\end{gather*}
$$

where we view ( ${ }^{t} a, a^{*}$ ) as an element of $\prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C})$ according to our convention of $\S 5.1$. Given $\Gamma$ as above, we can find a Z-lattice $M$ in $S$ and a subgroup $U$ of $G L_{n}(\mathfrak{r})$ of finite index such that $\left[\begin{array}{ll}1 & \sigma \\ 0 & 1\end{array}\right] \in \Gamma$ for every $\sigma \in M$ and $\operatorname{diag}[a, \widehat{a}] \in \Gamma$ for every $a \in U$. Thus, if $f \in \mathcal{M}_{\omega}(\Gamma)$, then

$$
\begin{align*}
& f(z+\sigma)=f(z) \quad \text { for every } \quad \sigma \in M  \tag{5.20}\\
& f\left(a z a^{*}\right)=\omega\left({ }^{t} a, a^{*}\right)^{-1} f(z) \quad \text { for every } \quad a \in U . \tag{5.21}
\end{align*}
$$

Notice that $\operatorname{tr}\left(\sigma \sigma^{\prime}\right) \in F$ for $\sigma, \sigma^{\prime} \in S$, and $\left(\sigma, \sigma^{\prime}\right) \mapsto \operatorname{Tr}_{F / \mathbf{Q}}\left(\operatorname{tr}\left(\sigma \sigma^{\prime}\right)\right)$ defines a nondegenerate pairing $S \times S \rightarrow \mathbf{Q}$. Let $L=\left\{h \in S \mid \operatorname{Tr}_{F / \mathbf{Q}}(\operatorname{tr}(h M)) \subset \mathbf{Z}\right\}$. Then (5.20) guarantees an expansion of the form

$$
\begin{equation*}
f(z)=\sum_{h \in L} c(h) \mathbf{e}_{\mathbf{a}}^{n}(h z) \tag{5.22a}
\end{equation*}
$$

with $c(h) \in X$. (For the proof, see [S97, Lemma A1.4].) We shall often put

$$
\begin{equation*}
f(z)=\sum_{h \in S} c(h) \mathbf{e}_{\mathbf{a}}^{n}(h z) \tag{5.22b}
\end{equation*}
$$

by defining $c(h)$ to be 0 for $h \in S, \notin L$. Usually we call the right-hand side of (5.22a) or (5.22b) the Fourier expansion of $f$, and call the $c(h)$ the Fourier coefficients of $f$. If $n=1$ and $F=\mathbf{Q}$, we take $\omega(x)=\operatorname{det}(x)^{k}$ with $k \in \mathbf{Z}$. Then (5.10) means: (5.23) For every $\alpha \in S L_{2}(\mathbf{Q})$ we have

$$
\left(f \|_{k} \alpha\right)(z)=\sum_{m=0}^{\infty} c_{\alpha m} \mathbf{e}\left(m z / N_{\alpha}\right)
$$

with $c_{\alpha m} \in \mathbf{C}$ and a positive integer $N_{\alpha}$.
In this situation we say that $f$ satisfies the cusp condition. Now if $n>1$ or $F \neq \mathbf{Q}$, no condition of this nature is necessary, because of the following fact:
5.7. Proposition. Suppose $n>1$ or $F \neq \mathbf{Q}$; let $f$ be a holomorphic function on $\mathcal{H}$ of the form (5.22a) satisfying (5.21) with a subgroup $U$ of $G L_{n}(\mathfrak{r})$ of finite index. Then $c(h) \neq 0$ only if $h_{v}$ is nonnegative for every $v \in \mathbf{a}$.

Proof. First we observe that $f(x+i y)=\sum_{h \in L} c(h) \mathbf{e}_{\mathbf{a}}^{n}(i h y) \mathbf{e}_{\mathbf{a}}^{n}(h x)$, and hence

$$
\begin{equation*}
\mathbf{e}_{\mathbf{a}}^{n}(i h y) c(h)=A \int_{S_{\mathbf{a}} / M} f(x+i y) \mathbf{e}_{\mathbf{a}}^{n}(-h x) d x \tag{5.24}
\end{equation*}
$$

where $A=\operatorname{vol}\left(S_{\mathbf{a}} / M\right)^{-1}$. Taking $X=\mathbf{C}^{t}$ with some $t$, put $\|w\|=\left(\sum_{k=1}^{t}\left|w_{k}\right|^{2}\right)^{1 / 2}$ for $w \in X$, and put also $\|\alpha\|=\operatorname{Max}_{\|w\|=1}\|\alpha w\|$ for $\alpha \in \operatorname{End}(X, \mathbf{C})$. Taking $y_{v}=$ $(2 \pi)^{-1} 1_{n}$ for every $v \in \mathbf{a}$ in (5.24), we obtain $\|c(h)\| \leq B \exp \left(\sum_{v \in \mathbf{a}} \operatorname{tr}(h)_{v}\right)$ with a constant $B$ independent of $h$. Now from (5.21) we obtain $c(h)=\omega\left({ }^{t} a, a^{*}\right)^{-1} c\left(a^{*} h a\right)$ for every $a \in U$, and hence

$$
\begin{equation*}
\|c(h)\| \leq B\left\|\omega\left({ }^{t} a, a^{*}\right)^{-1}\right\| \exp \left(\sum_{v \in \mathbf{a}} \operatorname{tr}\left(a^{*} h a\right)_{v}\right) \text { for every } a \in U \tag{*}
\end{equation*}
$$

Now suppose $n>1$; let $h$ be an element of $L$ such that $h_{u}$ is not nonnegative for some $u \in \mathbf{a}$. We consider Case UT; Case SP can be handled with obvious modifications. We can find $x \in\left(\mathbf{C}^{n}\right)^{\mathbf{a}}$ such that

$$
\begin{equation*}
\left(x_{1} x_{1}^{\rho}+x_{2} x_{2}^{\rho}\right)_{u}\left(x^{*} h x\right)_{u}<-\sum_{u \neq v \in \mathbf{a}}\left|\left(x_{1} x_{1}^{\rho}+x_{2} x_{2}^{\rho}\right)_{v}\left(x^{*} h x\right)_{v}\right| \tag{}
\end{equation*}
$$

where $x_{1}$ and $x_{2}$ are the first two components of $x$. (They are elements of $\mathbf{C}^{\mathbf{a}}$.) Since $K$ is dense in $\mathbf{C}^{\mathbf{a}}$, we can take such an $x$ in $K^{n}$. Multiplying it by a positive integer, we may assume that $x \in \mathfrak{r}^{n}$. Let $y=\left[\begin{array}{llll}-x_{2} & x_{1} & 0 & \cdots\end{array}\right]$ and $b=x y$; here $y \in\left(\mathbf{C}_{n}^{1}\right)^{\mathbf{a}}$ and $b \in\left(\mathbf{C}_{n}^{n}\right)^{\mathbf{a}}$. Then $b \in \mathfrak{r}_{n}^{n}$ and $b^{2}=0$ since $y x=0$. We have
(***)

$$
\operatorname{tr}\left(b^{*} h b\right)=\operatorname{tr}\left(y^{*} x^{*} h x y\right)=y y^{*} x^{*} h x=\left(x_{1} x_{1}^{\rho}+x_{2} x_{2}^{\rho}\right)\left(x^{*} h x\right)
$$

Put $a=(1+b)^{m}$ with $0<m \in \mathbf{Z}$. Since $1+b \in S L_{n}(\mathfrak{r})$, we see that $a \in U$ if $m$ is a multiple of some positive integer $N$. We have

$$
\operatorname{tr}\left(a^{*} h a\right)=\operatorname{tr}\left((1+m b)^{*} h(1+m b)\right)=\operatorname{tr}(h)+m \cdot \operatorname{tr}\left(h\left(b+b^{*}\right)\right)+m^{2} \cdot \operatorname{tr}\left(b^{*} h b\right) .
$$

Thus $\sum_{v \in \mathbf{a}} \operatorname{tr}\left(a^{*} h a\right)_{v}=p+m q+m^{2} r$ with $p, q, r \in \mathbf{R}$, and $r=\sum_{v \in \mathbf{a}}\left(x_{1} x_{1}^{\rho}+\right.$ $\left.x_{2} x_{2}^{\rho}\right)_{v}\left(x^{*} h x\right)_{v}<0$ by $\left({ }^{* * *}\right)$ and $\left({ }^{* *}\right)$. Now by $\left({ }^{*}\right)$,

$$
\|c(h)\| \leq B\left\|\omega\left(1+{ }^{t} b, 1+b^{*}\right)^{-1}\right\|^{m} \exp \left(p+m q+m^{2} r\right)
$$

for $0<m \in N \mathbf{Z}$. Making $m$ large, we find that $c(h)=0$ as expected.
Next suppose $n=1$ and $F \neq \mathbf{Q}$. In this case we have $S=F$ and we can take $U$ to be a subgroup of $\mathfrak{g}^{\times}$of finite index. Let $h$ be an element of $F$ such that $h_{u}<0$ for some $u \in \mathbf{a}$. Since $F \neq \mathbf{Q}$, we can find an element $a \in U$ such that $\left|a_{u}\right|>1$ and $\left|a_{v}\right|<1$ for $u \neq v \in \mathbf{a}$. Now $\operatorname{tr}\left(\left(a^{m}\right)^{*} h a^{m}\right)=h a^{2 m}$ in this case, and therefore we obtain $\|c(h)\|=0$ by taking $a^{m}$ in place of $a$ in $\left(^{*}\right)$ and making $m$ large.

Another type of proof of Proposition 5.7 for the elements of $\mathcal{M}_{k}, k \in \mathbf{Z}^{\mathbf{a}}$, is given in [S97, Proposition A4.2]. See also [S78b, Proposition 3.1] and [S97, Proposition A4.5] for some results of the same nature in different settings.
5.8. Let $f \in \mathcal{M}_{\omega}(\Gamma)$ and $\alpha \in \widetilde{G}_{+}$. Then we easily see that $f \|_{\omega} \alpha \in \mathcal{M}_{\omega}\left(\alpha^{-1} \Gamma \alpha\right)$. Thus $\mathcal{M}_{\omega}$ is stable under the map $f \mapsto f \|_{\omega} \alpha$ for every $\alpha \in \widetilde{G}_{+}$. Now for $f \in$ $\mathcal{M}_{\omega}(\Gamma)$ and $\alpha \in \widetilde{G}_{+}$we have an expansion

$$
\begin{equation*}
\left(f \|_{\omega} \alpha\right)(z)=\sum_{h \in S} c_{\alpha}(h) \mathbf{e}_{\mathbf{a}}^{n}(h z) \tag{5.25}
\end{equation*}
$$

By Proposition 5.7, $c_{\alpha}(h) \neq 0$ only if $h_{v}$ is nonnegative for every $v \in \mathbf{a}$. We call $f$ a cusp form if $c_{\alpha}(h)=0$ for every $\alpha \in G$ and for every $h$ such that $\operatorname{det}(h)=0$, and denote by $\mathcal{S}_{\omega}(\Gamma)$ (resp. $\mathcal{S}_{\omega}$ ) the set of all cusp forms contained in $\mathcal{M}_{\omega}(\Gamma)$ (resp. $\mathcal{M}_{\omega}$ ). In view of Lemma 1.3 (3), if $f$ is a cusp form, then $c_{\alpha}(h)=0$ for every $\alpha \in \widetilde{G}_{+}$and for every $h$ such that $\operatorname{det}(h)=0$. If $f \in \mathcal{S}_{\omega}(\Gamma)$ and $\alpha \in \widetilde{G}_{+}$, then $f \|_{\omega} \alpha \in \mathcal{S}_{\omega}\left(\alpha^{-1} \Gamma \alpha\right)$.

To consider the arithmeticity of modular forms, let us hereafter assume that $(X, \omega)$ has a $\mathbf{Q}$-structure in the sense that $X=X_{0} \otimes_{\mathbf{Q}} \mathbf{C}$ with a fixed vector space $X_{0}$ over $\mathbf{Q}$, and $\omega$ is the natural extension of a rational representation $\omega_{0}$ : $G L_{n}(\mathbf{Q})^{\mathbf{b}} \rightarrow G L\left(X_{0}\right)$. (Often $X=\mathbf{C}^{m}, X_{0}=\mathbf{Q}^{m}$, and $\omega_{0}$ is a representation $\left.G L_{n}(\mathbf{Q})^{\mathbf{b}} \rightarrow G L_{m}(\mathbf{Q}).\right)$ Then, given a subfield $D$ of $\mathbf{C}$, we say that $f$ of (5.23) is $D$-rational if $c(h) \in X_{0} \otimes_{\mathbf{Q}} D$ for all $h \in S$, and denote by $\mathcal{M}_{\omega}(D)$ the set of all $D$-rational elements of $\mathcal{M}_{\omega}$. Then we put

$$
\begin{align*}
\mathcal{A}_{\omega}(D) & =\bigcup_{e}\left\{p^{-1} q \mid q \in \mathcal{M}_{\tau_{e}}(D), 0 \neq p \in \mathcal{M}_{e}(D)\right\}  \tag{5.26a}\\
\mathcal{M}_{\omega}(\Gamma, D) & =\mathcal{M}_{\omega}(\Gamma) \cap \mathcal{M}_{\omega}(D), \quad \mathcal{A}_{\omega}(\Gamma, D)=\mathcal{A}_{\omega}(\Gamma) \cap \mathcal{A}_{\omega}(D)  \tag{5.26b}\\
\mathcal{S}_{\omega}(\Gamma, D) & =\mathcal{S}_{\omega}(\Gamma) \cap \mathcal{M}_{\omega}(D), \quad \mathcal{S}_{\omega}(D)=\mathcal{S}_{\omega} \cap \mathcal{M}_{\omega}(D) \tag{5.26c}
\end{align*}
$$

where $\tau_{e}(x)=\operatorname{det}(x)^{e} \omega(x), e \in \mathbf{Z}^{\mathbf{b}}$. We use the subscript $k$ instead of $\omega$ (that is, we write $\mathcal{A}_{k}(D)$ and $\mathcal{S}_{k}$ for $\mathcal{A}_{\omega}(D)$ and $\mathcal{S}_{\omega}$, for example) if $\omega(x)=\operatorname{det}(x)^{k}$.

Clearly $\mathcal{A}_{0}(D)$ (resp. $\mathcal{A}_{0}(\Gamma, D)$ ) is a subfield of $\mathcal{A}_{0}$ (resp. $\mathcal{A}_{0}(\Gamma)$ ). Also $E \mathcal{A}_{0}(\Gamma, D) \subset \mathcal{A}_{0}(\Gamma, E)$ (resp. $E \mathcal{A}_{0}(D) \subset \mathcal{A}_{0}(E)$ ) if $E$ is an extension of $D$. The equalities $E \mathcal{A}_{0}(\Gamma, D)=\mathcal{A}_{0}(\Gamma, E)$ and $E \mathcal{A}_{0}(D)=\mathcal{A}_{0}(E)$ are true in some cases, but they are not necessarily true in general. It should also be noted that $\mathcal{M}_{k}$ can be $\{0\}$ even if $k_{v}>0$ for every $v \in \mathbf{a} ; \mathcal{S}_{k}$ can be $\{0\}$ even if $\mathcal{M}_{k} \neq\{0\}$; see Proposition 6.16 below.
5.9. For $h \in S$ we write $0 \leq h$ or $h \geq 0$ if $h_{v}$ is nonnegative for every $v \in \mathbf{a}$. By Proposition 5.7 the expansion of $f$ in (5.22a) can be written

$$
\begin{equation*}
f(z)=\sum_{0 \leq h \in L} c(h) \mathbf{e}_{\mathbf{a}}^{n}(h z) \tag{5.27}
\end{equation*}
$$

Let $\mathfrak{T}(L)$ denote the set of all formal series of the form (5.27) with $c(h) \in \mathbf{C}$, and $\mathfrak{T}$ the union of $\mathfrak{T}(L)$ for all Z-lattices $L$ in $S$. Then $\mathfrak{T}$ has a natural ring-structure. Put $[x]=\operatorname{Tr}_{F / \mathbf{Q}}(\operatorname{tr}(x))$ for $x \in S$ and let $m$ be the dimension of $S$ over $\mathbf{Q}$. Then we can find a Q-basis $\left\{s_{1}, \ldots, s_{m}\right\}$ of $S$ such that $\left[s_{i} L\right] \subset \mathbf{Z}$ and $s_{i v}>0$ for every $i$ and every $v \in \mathbf{a}$. Clearly $\left[s_{i} h\right] \geq 0$ if $0 \leq h \in L$. Taking $m$ independent indeterminates $\xi_{1}, \ldots, \xi_{m}$, for $f \in \mathfrak{T}(L)$ as in (5.27) we put

$$
\begin{equation*}
\psi(f)=\sum_{0 \leq h \in L} c(h) \prod_{i=1}^{m} \xi_{i}^{\left[s_{i} h\right]} . \tag{5.28}
\end{equation*}
$$

We easily see that $\psi$ defines a ring-injection of $\mathfrak{T}(L)$ into the ring $\mathbf{C}\left[\left[\xi_{1}, \ldots, \xi_{m}\right]\right]$ of all formal power series in $\xi_{1}, \ldots, \xi_{m}$ with coefficients in C. Therefore $\mathfrak{T}(L)$ is an integral domain, and the same is true for $\mathfrak{T}$. Given $\sigma \in \operatorname{Aut}(\mathbf{C})$, we obtain automorphisms of $\mathbf{C}\left[\left[\xi_{1}, \ldots, \xi_{m}\right]\right]$ and $\mathfrak{T}$ by applying $\sigma$ to the coefficients; denote by $f^{\sigma}$ the image of $f$ under these automorphisms. This means that for $f \in \mathfrak{T}(L)$ as in (5.27) we have

$$
\begin{equation*}
f^{\sigma}=\sum_{0 \leq h \in L} c(h)^{\sigma} \mathbf{e}_{\mathbf{a}}^{n}(h z) \tag{5.29}
\end{equation*}
$$

and clearly $\psi(f)^{\sigma}=\psi\left(f^{\sigma}\right)$. We can in fact define $f^{\sigma}$ formally in the same manner even when $c(h) \in X$, since $\sigma$ acts naturally on $X$. This action of $\sigma$ can be extended to $\mathcal{A}_{\omega}$. Indeed, for $f=p^{-1} q \in \mathcal{A}_{\omega}$ with $0 \neq p \in \mathcal{M}_{e}$ and $q \in \mathcal{M}_{\tau_{e}}$ as in (5.26a) we put $f^{\sigma}=\left(p^{\sigma}\right)^{-1} q^{\sigma}$. This is a vector whose components belong to the field of quotients of $\mathfrak{T}$. Clearly this is well-defined. For the moment, $f^{\sigma}$ is merely defined formally, and, in general, not defined as a function on $\mathcal{H}$. However, we shall later show that it is always meaningful as a function on $\mathcal{H}$.

Let us now prove

$$
\begin{equation*}
\mathcal{A}_{\omega}(D) \cap \mathcal{M}_{\omega}=\mathcal{M}_{\omega}(D) \tag{5.30}
\end{equation*}
$$

If $f \in \mathcal{A}_{\omega}(D) \cap \mathcal{M}_{\omega}$, then $f=p^{-1} q$ with $0 \neq p \in \mathcal{M}_{e}(D)$ and $q \in \mathcal{M}_{\tau_{e}}(D)$. Then for $\sigma \in \operatorname{Aut}(\mathbf{C} / D)$ we have $p f^{\sigma}=(p f)^{\sigma}=q^{\sigma}=q=p f$. Since $p$ is not a zero-divisor, we obtain $f^{\sigma}=f$, that is, $f \in \mathcal{M}_{\omega}(D)$.

There are several natural, but highly nontrivial, questions concerning $\mathcal{M}_{\omega}$ :
(Q1) If $f \in \mathcal{M}_{\omega}$ and $\sigma \in \operatorname{Aut}(\mathbf{C})$, does the formal series for $f^{\sigma}$ define an automorphic form? If so, what is its weight?
(Q2) Can we find an algebraic number field $D$ such that $\mathcal{M}_{\omega}=\mathcal{M}_{\omega}(D) \otimes_{D} \mathbf{C}$ ? When can we take $D=\mathbf{Q}$ ?
(Q3) Given $\alpha \in \widetilde{G}_{+}$and $f \in \mathcal{M}_{\omega}$, how is $\left(f \|_{\omega} \alpha\right)^{\sigma}$ related.to $f^{\sigma}$ ? Can we find $\beta \in \widetilde{G}_{+}$and a weight $\psi$ such that $\left(f \|_{\omega} \alpha\right)^{\sigma}=f^{\sigma} \|_{\psi} \beta$ ?

We can ask similar questions on $\mathcal{S}_{\omega}$ and $\mathcal{A}_{\omega}$. We shall answer these questions in Sections 9 and 10. In connection with (Q2) we note here an easy lemma.
5.10. Lemma. (1) If $D$ is a subfield of $\mathbf{C}$ and $f_{1}, \ldots, f_{m}$ are elements of $\mathcal{M}_{\omega}(D)$ linearly independent over $D$, then they are linearly independent over $\mathbf{C}$.
(2) If $\mathcal{M}_{\omega}(\Gamma, D)$ spans $\mathcal{M}_{\omega}(\Gamma)$ over $\mathbf{C}$, then $\mathcal{M}_{\omega}(\Gamma)=\mathcal{M}_{\omega}(\Gamma, D) \otimes_{D} \mathbf{C}$.

Proof. Let $f_{i}(z)=\sum_{h} c_{i}(h) \mathbf{e}_{\mathbf{a}}^{n}(h z)$. Put $W_{h}=\left\{x \in \mathbf{C}^{m} \mid \sum_{i=1}^{m} x_{i} c_{i}(h)=0\right\}$ for $h \in S$ and $Y=\bigcap_{h \in S} W_{h}$. Then each $W_{h}$, as well as $Y$, is a vector subspace of $\mathbf{C}^{m}$ defined over $D$. Since $Y$ has no $D$-rational point other than 0 , we have $Y=\{0\}$, which proves (1). Then (2) follows immediately from (1).
5.11. There is one phenomenon peculiar to Case UT. First put ${ }^{t} z=\left({ }^{t} z_{v}\right)_{v \in \mathbf{a}}$ for $z=\left(z_{v}\right)_{v \in \mathbf{a}} \in \mathcal{H}$. Let $\alpha \in \widetilde{G}_{+}$. Then $\alpha^{\rho} \in \widetilde{G}_{+}$and from (3.14), (3.15), and (3.16) we can easily derive that

$$
\begin{equation*}
\alpha^{\rho}\left({ }^{t} z\right)={ }^{t}(\alpha z), \quad \lambda(\alpha, z)=\mu\left(\alpha^{\rho},{ }^{t} z\right) \tag{5.31}
\end{equation*}
$$

Given $\{\omega, X\}$ and $f: \mathcal{H} \rightarrow X$, define $f^{\prime}: \mathcal{H} \rightarrow X$ by $f^{\prime}(z)=f\left({ }^{t} z\right)$. Then we easily see that

$$
\begin{equation*}
\left(f \|_{\omega} \alpha\right)^{\prime}=f^{\prime} \|_{\omega^{\rho}} \alpha^{\rho} \tag{5.32}
\end{equation*}
$$

where $\omega^{\rho}$ is defined by $\omega^{\rho}(a, b)=\omega(b, a)$. Therefore $f^{\prime} \in \mathcal{M}_{\omega^{\rho}}$ if $f \in \mathcal{M}_{\omega}$.
Define also $\widetilde{f}=f^{\sim}: \mathcal{H} \rightarrow X$ by $\widetilde{f}(z)=\overline{f\left(-z^{*}\right)}$ and put $\alpha^{\prime}=\varepsilon \alpha \varepsilon$ for $\alpha \in \widetilde{G}_{+}$ with $\varepsilon=\operatorname{diag}\left[1_{r},-1_{r}\right]$. Then we can easily verify that

$$
\begin{equation*}
\alpha\left(-z^{*}\right)=-\left(\alpha^{\prime} z\right)^{*}, \quad \lambda\left(\alpha,-z^{*}\right)=\overline{\mu\left(\alpha^{\prime}, z\right)}, \tag{5.33}
\end{equation*}
$$

$$
\begin{equation*}
\left(f \|_{\omega} \alpha\right)^{\sim}=\widetilde{f} \|_{\omega^{\rho}} \alpha^{\prime} \tag{5.34}
\end{equation*}
$$

provided $\omega$ is R-rational, in which case $\tilde{f} \in \mathcal{M}_{\omega^{\rho}}$ if $f \in \mathcal{M}_{\omega}$.
5.12. The measure $\mathbf{d} z$ of Lemma 3.4 in Case SP can be written

$$
\begin{equation*}
\mathbf{d}(x+i y)=\operatorname{det}(y)^{-n-1} d x d y \quad(\text { Case SP }) \tag{5.35}
\end{equation*}
$$

with $d x=\prod_{h \leq k} d x_{h k}$ and $d y=\prod_{h \leq k} d y_{h k}$ for real symmetric matrices $x$ and $y$. In Case UT, for each fixed $v \in$ a we have $\mathbf{C}_{n}^{n}=S_{v} \oplus i S_{v}$ with $S$ of (5.17). Now $S_{v}$ is the vector space of all hermitian matrices of size $n$. We can identify $S_{v}$ with $\mathbf{R}^{n^{2}}$ through the map $x \mapsto\left(x_{h h}, \operatorname{Re}\left(x_{h k}\right), \operatorname{Im}\left(x_{h k}\right)(h<k)\right)$, and define a measure $d x$ on $S_{v}$ by pulling back the standard measure on $\mathbf{R}^{n^{2}}$. Writing $z=\left(z_{h k}\right) \in \mathbf{C}_{n}^{n}$ in the form $z=x+i y$ with $x, y \in S_{v}$, we have a measure $d x d y$ on $\mathbf{C}_{n}^{n}$ with the measures $d x$ and $d y$ on $S_{v}$ given as above. It should be noted that this does not coincide with the standard measure $\prod_{h, k}\left[(i / 2) d z_{h k} \wedge d \bar{z}_{h k}\right]$ on $\mathbf{C}_{n}^{n}$. In fact, we easily see that

$$
\begin{equation*}
\prod_{h, k}\left[(i / 2) d z_{h k} \wedge d \bar{z}_{h k}\right]=2^{n(n-1)} d x d y \tag{5.36}
\end{equation*}
$$

Thus the measure $\mathbf{d} z$ of Lemma 3.4 in Case UT can be given by

$$
\begin{equation*}
\mathbf{d}(x+i y)=2^{n(n-1)} \operatorname{det}(y)^{-2 n} d x d y \quad\left(x^{*}=x, y^{*}=y ; \text { Case UT }\right) \tag{5.37}
\end{equation*}
$$

## 6. Parametrization by theta functions

6.1. In this section we introduce certain theta functions, by which we parametrize our abelian varieties in Cases SP and UT. If the basic field is $\mathbf{Q}$, they are the classical $\theta$ and its modification $\varphi$ given by

$$
\begin{align*}
& \theta(u, z ; r, s)=\sum_{g-r \in \mathbf{Z}^{n}} \mathbf{e}\left(2^{-1} \cdot{ }^{t} g z g+{ }^{t} g(u+s)\right),  \tag{6.1}\\
& \varphi(u, z ; r, s)=\mathbf{e}\left(2^{-1} \cdot{ }^{t} u(z-\bar{z})^{-1} u\right) \theta(u, z ; r, s) . \tag{6.2}
\end{align*}
$$

Here $u \in \mathbf{C}^{n}, z \in \mathfrak{H}_{n}, r, s \in \mathbf{R}^{n}$, and $\mathbf{e}(c)=\exp (2 \pi i c)$ as we set in (5.14). We are going to consider the pullbacks of these (with larger $n$ ) to $\left(\mathbf{C}^{n}\right)^{\mathbf{b}} \times \mathcal{H}$, which may be viewed also as generalizations of (6.1) and (6.2).

We need some new symbols. First, we denote by $\mathbf{Q}_{\mathrm{ab}}$ the maximal abelian extension of $\mathbf{Q}$ in $\mathbf{C}$. Next, given a finite-dimensional vector space $W$ over $\mathbf{Q}$, we denote by $\mathcal{S}\left(W_{\mathbf{h}}\right)$ the Schwartz-Bruhat space of $W_{\mathbf{h}}$ (see $\S 1.6$ ). Also, for a square matrix $x$ of size $2 n$ we denote by $a_{x}, b_{x}, c_{x}$, and $d_{x}$ the $n \times n$-blocks of $x$ in the sense of $\S 1.8$.

Let $[F: \mathbf{Q}]=e$. In this section we often identify a with $\{1, \ldots, e\}$ so that $X^{\mathbf{a}}$ may be written $X^{e}$ for various symbols $X$. For example, $\left(\mathbf{C}^{n}\right)^{\mathbf{a}}$ can be identified with $\mathbf{C}^{e n}$, and $\operatorname{diag}\left[a_{v}\right]_{v \in \mathbf{a}}$ with $\operatorname{diag}\left[a_{1}, \ldots, a_{e}\right]$.
6.2. Let $\left\{\beta_{1}, \ldots, \beta_{e}\right\}$ be a $\mathbf{Q}$-basis of $F$, and $\left\{\gamma_{1}, \ldots, \gamma_{e}\right\}$ be another $\mathbf{Q}$ basis of $F$ determined by the condition $\operatorname{Tr}_{F / \mathbf{Q}}\left(\beta_{i} \gamma_{j}\right)=\delta_{i j}$. Define a $\mathbf{Q}$-linear map $g: \mathbf{Q}_{2 e n}^{1} \rightarrow F_{2 n}^{1}$ by

$$
\begin{equation*}
g\left(x_{1}, \ldots, x_{e}, y_{1}, \ldots, y_{e}\right)=\left(\sum_{i=1}^{e} \beta_{i} x_{i}, \sum_{i=1}^{e} \gamma_{i} y_{i}\right) \quad\left(x_{i}, y_{i} \in \mathbf{Q}_{n}^{1}\right) \tag{6.3}
\end{equation*}
$$

A simple calculation shows that

$$
\begin{equation*}
\operatorname{Tr}_{F / \mathbf{Q}}\left(g(u) \eta_{n} \cdot{ }^{t} g\left(u^{\prime}\right)\right)=u \eta_{e n} \cdot{ }^{t} u^{\prime} \tag{6.4}
\end{equation*}
$$

for $u, u^{\prime} \in \mathbf{Q}_{2 e n}^{1}$. Denoting the $j$-th conjugate of $\beta_{i}$ by $\beta_{i j}$, put

$$
\begin{gather*}
B=\left[\begin{array}{ccc}
\beta_{11} 1_{n} & \cdots & \beta_{1 e} 1_{n} \\
\cdots & \cdots & \cdots \\
\beta_{e 1} 1_{n} & \cdots & \beta_{e e} 1_{n}
\end{array}\right],  \tag{6.5}\\
\left.\psi(a)=\operatorname{diag}\left[a_{v}\right]_{v \in \mathbf{a}}=\operatorname{diag}\left[a_{1}, \ldots, a_{e}\right] \quad\left(\left(a_{v}\right)_{v \in \mathbf{a}} \in\left(\mathbf{C}_{n}^{n}\right)^{\mathbf{a}}=\left(\mathbf{C}_{n}^{n}\right)^{e}\right)\right),  \tag{6.6}\\
\omega(\alpha)=\left[\begin{array}{cc}
B & 0 \\
0 & { }^{t} B^{-1}
\end{array}\right]\left[\begin{array}{cc}
\psi\left(a_{\alpha}\right) & \psi\left(b_{\alpha}\right) \\
\psi\left(c_{\alpha}\right) & \psi\left(d_{\alpha}\right)
\end{array}\right]\left[\begin{array}{cc}
B^{-1} & 0 \\
0 & { }^{t} B
\end{array}\right] \quad\left(\alpha \in F_{2 n}^{2 n}\right) . \tag{6.7}
\end{gather*}
$$

Then we can easily verify that $\omega(\alpha) \in \mathbf{Q}_{2 e n}^{2 e n}$ and

$$
\begin{equation*}
g(u \omega(\alpha))=g(u) \alpha \quad\left(u \in \mathbf{Q}_{2 e n}^{1}, \alpha \in F_{2 n}^{2 n}\right) . \tag{6.8}
\end{equation*}
$$

Let us now define a subgroup $G_{0}=G_{0}\left(F, \eta_{n}\right)$ of $G p(n, F)$ by

$$
\begin{equation*}
G_{0}=G_{0}\left(F, \eta_{n}\right)=\{\alpha \in G p(n, F) \mid \nu(\alpha) \in \mathbf{Q}\} . \tag{6.9}
\end{equation*}
$$

This is the same as the group of (3.29). In view of (6.4) we can show that $\omega$ defines an injective homomorphism of $G_{0}\left(F, \eta_{n}\right)$ into $G p(e n, \mathbf{Q})$, and also an injection of $S p(n, F)$ into $S p(e n, \mathbf{Q})$. Define also an embedding $\varepsilon: \mathfrak{H}_{n}^{\mathrm{a}} \rightarrow \mathfrak{H}_{e n}$ by

$$
\begin{equation*}
\varepsilon(z)=B \cdot \operatorname{diag}\left[z_{v}\right]_{v \in \mathbf{a}} \cdot{ }^{t} B \quad\left(z=\left(z_{v}\right)_{v \in \mathbf{a}} \in \mathfrak{H}_{n}^{\mathbf{a}}\right) . \tag{6.10}
\end{equation*}
$$

Then we can easily verify that

$$
\begin{align*}
& \omega(\alpha) \varepsilon(z)=\varepsilon(\alpha z) \quad\left(\alpha \in G_{0}\left(F, \eta_{n}\right), \nu(\alpha)>0\right)  \tag{6.11a}\\
& \mu(\omega(\alpha), \varepsilon(z))={ }^{t} B^{-1} \operatorname{diag}\left[\mu_{v}(\alpha, z)\right]_{v \in \mathbf{a}} \cdot{ }^{t} B \tag{6.11b}
\end{align*}
$$

6.3. Let us describe the above embedding $\varepsilon$ in terms of the families of abelian varieties established in $\S 4.7$. We consider two families

$$
\begin{aligned}
\mathcal{F}\left(\Omega_{0}\right)= & \left\{\mathcal{P}_{z} \mid z \in \mathfrak{H}_{e n}\right\}, \quad \mathcal{P}_{z}=\left(A_{z}, \mathcal{C}_{z}, \iota_{z} ;\left\{t_{i}(z)\right\}_{i=1}^{s}\right) \\
\mathcal{F}\left(\Omega_{F}\right)= & \left\{\mathcal{P}_{z} \mid z \in \mathfrak{H}_{n}^{\mathrm{a}}\right\}, \quad \mathcal{P}_{z}=\left(A_{z}, \mathcal{C}_{z}, \iota_{z} ;\left\{t_{i}(z)\right\}_{i=1}^{s}\right) \\
& \Omega_{0}=\left\{\mathbf{Q}, \text { id., } L_{0}, \eta_{e n},\left\{u_{i}\right\}_{i=1}^{s}\right\}, \\
& \Omega_{F}=\left\{F, \Psi, L_{F}, \eta_{n},\left\{g\left(u_{i}\right)\right\}_{i=1}^{s}\right\} .
\end{aligned}
$$

Here $L_{0}$ resp. $L_{F}$ is a $\mathbf{Z}$-lattice in $\mathbf{Q}_{2 e n}^{1}$ resp. $F_{n}^{1}$. By (4.15) we have to assume that $\operatorname{Tr}_{F / \mathbf{Q}}\left(x \eta_{n} \cdot{ }^{t} y\right) \in \mathbf{Z}$ for every $x, y \in L_{F}$. We take $L_{0}=g^{-1}\left(L_{F}\right)$. By (6.4), $u \eta_{e n} \cdot{ }^{t} u^{\prime} \in \mathbf{Z}$ for every $u, u^{\prime} \in L_{0}$. Thus $\Omega_{0}$ is meaningful for this $L_{0}$. Now recall that for each $z \in \mathfrak{H}_{n}^{\text {a }}$ the variety $A_{z}$ can be given by $\left(\mathbf{C}^{n}\right)^{\mathbf{a}} / p_{z}\left(L_{F}\right)$ with the map $p_{z}$ of (4.24). Disregarding the endomorphism algebra $\iota_{z}(F)$, we observe that $\mathcal{P}_{z}$ is of type $\Omega_{0}$, and so isomorphic to $\mathcal{P}_{w}$ with some $w \in \mathfrak{H}_{e n}$. We obtain such a $w$ by checking the period matrix $X\left(p_{z} \circ g\right)$ of (4.18). In fact, a simple calculation shows that the upper half of $X\left(p_{z} \circ g\right)$ is of the form

$$
\left[\begin{array}{cccccc}
\beta_{11} z_{1} & \cdots & \beta_{e 1} z_{1} & \gamma_{11} 1_{n} & \cdots & \gamma_{e 1} 1_{n} \\
\cdots & \cdots & \cdots & \cdots & \cdots & \cdots \\
\beta_{1 e} z_{e} & \cdots & \beta_{e e} z_{e} & \gamma_{1 e} 1_{n} & \cdots & \gamma_{e e} 1_{n}
\end{array}\right] .
$$

This equals $B^{-1}\left[\varepsilon(z) \quad 1_{e n}\right]$. Since $p_{z}(g(x))$ equals this matrix times ${ }^{t} x$ for $x \in$ $\mathbf{Q}_{2 \text { en }}^{1}$, we obtain

$$
\begin{equation*}
B \cdot p_{z}(g(x))=p_{\varepsilon(z)}(x) \quad\left(x \in \mathbf{Q}_{2 e n}^{1}\right) . \tag{6.12}
\end{equation*}
$$

Therefore the map $u \mapsto B u$ for $u \in\left(\mathbf{C}^{n}\right)^{\mathbf{a}}$ gives an isomorphism of $\left(\mathbf{C}^{n}\right)^{\mathbf{a}} / p_{z}\left(L_{F}\right)$ onto $\mathbf{C}^{e n} / p_{\varepsilon(z)}\left(L_{0}\right)$, or rather, an isomorphism of $\mathcal{P}_{z}$ (minus $\iota_{z}$ ) to $\mathcal{P}_{\varepsilon(z)}$ because of (6.4).
6.4. We now define theta functions $\theta_{F}$ and $\varphi_{F}$ by

$$
\begin{align*}
& \theta_{F}(u, z ; r, s)=\sum_{h-r \in M} \mathbf{e}_{\mathbf{a}}\left(2^{-1} \cdot{ }^{t} h z h+{ }^{t} h(u+s)\right),  \tag{6.13}\\
& \varphi_{F}(u, z ; r, s)=\mathbf{e}_{\mathbf{a}}\left(2^{-1} \cdot{ }^{t} u(z-\bar{z})^{-1} u\right) \theta_{F}(u, z ; r, s) \tag{6.14}
\end{align*}
$$

Here $u \in\left(\mathbf{C}^{n}\right)^{\mathbf{a}}, z \in \mathfrak{H}_{n}^{\mathbf{a}}, r, s \in\left(\mathbf{R}^{n}\right)^{\mathbf{a}}$, and $\mathbf{e}_{\mathbf{a}}(x)=\exp \left(\sum_{v \in \mathbf{a}} x_{v}\right)$ for $x \in \mathbf{C}^{\mathbf{a}}$ as we set in (5.15); $M$ will be specified after (6.15). If $F=\mathbf{Q}$ and $M=\mathbf{Z}^{n}$, then these coincide with the functions of (6.1) and (6.2). By an easy calculation we can verify that these are the pullbacks of (6.1) and (6.2) in the sense that

$$
\begin{align*}
\varphi_{F}\left(u, z ; g_{1}(r), g_{2}(s)\right)=\varphi & (B u, \varepsilon(z) ; r, s)  \tag{6.15}\\
& \left(u \in\left(\mathbf{C}^{n}\right)^{\mathbf{a}}, z \in \mathfrak{H}_{n}^{\mathbf{a}}, r, s \in \mathbf{Q}^{2 e n}\right),
\end{align*}
$$

where $g_{1}(q)=\sum_{i=1}^{e} \beta_{i} q_{i}$ and $g_{2}(q)=\sum_{i=1}^{e} \gamma_{i} q_{i}$ for ${ }^{t} q=\left[\begin{array}{lll}{ }^{t} q_{1} & \ldots & { }^{t} q_{e}\end{array}\right]$ with $q_{i} \in$ $\mathbf{Q}^{n}$, and the same type of equality holds with $\theta$ in place of $\varphi$. Now $M=g_{1}\left(\mathbf{Z}^{e n}\right)$. It is not difficult to show that the right-hand side of (6.13) is locally uniformly convergent on $\left(\mathbf{C}^{n}\right)^{\mathbf{a}} \times\left(\mathfrak{H}_{n}\right)^{\mathbf{a}}$, and so defines a holomorphic function in $(u, z)$.

For our later purposes it is convenient to consider functions of the following types:

$$
\begin{align*}
& \theta_{F}(u, z ; \lambda)=\sum_{h \in F^{n}} \lambda(h) \mathbf{e}_{\mathbf{a}}\left(2^{-1} \cdot{ }^{t} h z h+{ }^{t} h u\right),  \tag{6.16}\\
& \varphi_{F}(u, z ; \lambda)=\mathbf{e}_{\mathbf{a}}\left(2^{-1} \cdot{ }^{t} u(z-\bar{z})^{-1} u\right) \theta_{F}(u, z ; \lambda) . \tag{6.17}
\end{align*}
$$

Here $\lambda \in \mathcal{S}\left(F_{\mathrm{h}}^{n}\right)$. Clearly (6.13) (resp. (6.14)) is a special case of (6.16) (resp. (6.17)). As we said in §1.6, we view $\lambda$ as a function on $F_{\mathbf{A}}^{n}$, so that $\lambda(h)$ for $h \in F^{n}$ is meaningful.
6.5. Let $K$ be a totally imaginary quadratic extension of $F$ as in $\S 3.5$. Take $\zeta \in K$ so that $\zeta^{\rho}=-\zeta$ and $K=F(\zeta)$. Define an $F$-linear map $h: F_{4 n}^{1} \rightarrow K_{2 n}^{1}$ by

$$
\begin{equation*}
h\left(x_{1}, x_{2}, y_{1}, y_{2}\right)=\left(x_{1}-\zeta x_{2}, 2^{-1} y_{1}+(2 \zeta)^{-1} y_{2}\right) \quad\left(x_{i}, y_{i} \in F_{n}^{1}\right) \tag{6.18}
\end{equation*}
$$

Then we have $\operatorname{Tr}_{K / F}\left(h(u) \eta_{n} h\left(u^{\prime}\right)^{*}\right)=u \eta_{2 n} \cdot{ }^{t} u^{\prime}$ for $u, u^{\prime} \in F_{4 n}^{1}$.
We can now define embeddings $\tau: K_{2 n}^{2 n} \rightarrow F_{4 n}^{4 n}$ and $\psi: \mathcal{H}_{n}^{\mathbf{a}} \rightarrow \mathfrak{H}_{2 n}^{\mathrm{a}}$ as follows:

$$
\begin{align*}
& A=\left[\begin{array}{cc}
1_{n} & 1_{n} \\
\zeta 1_{n} & -\zeta 1_{n}
\end{array}\right], \quad \sigma(a)=\operatorname{diag}\left[a^{\rho}, a\right] \quad\left(a \in K_{n}^{n}\right),  \tag{6.19}\\
& \tau(\alpha)=\left[\begin{array}{cc}
A & 0 \\
0 & \widehat{A}
\end{array}\right]\left[\begin{array}{ll}
\sigma\left(a_{\alpha}\right) & \sigma\left(b_{\alpha}\right) \\
\sigma\left(c_{\alpha}\right) & \sigma\left(d_{\alpha}\right)
\end{array}\right]\left[\begin{array}{cc}
A^{-1} & 0 \\
0 & A^{*}
\end{array}\right] \quad\left(\alpha \in K_{2 n}^{2 n}\right),  \tag{6.20}\\
& \psi(w)=\left(A_{v} \cdot \operatorname{diag}\left[{ }^{t} w_{v}, w_{v}\right] A_{v}^{*}\right)_{v \in \mathbf{a}} \quad\left(w=\left(w_{v}\right)_{v \in \mathbf{a}} \in \mathcal{H}_{n}^{\mathbf{a}}\right) . \tag{6.21}
\end{align*}
$$

Then we can easily verify that $\tau$ defines an injective homomorphism of $G U\left(\eta_{n}\right)$ into $G p(n, F)$ and

$$
\begin{array}{lc}
h(u \tau(\alpha))=h(u) \alpha & \left(u \in F_{4 n}^{1}, \alpha \in F_{4 n}^{4 n}\right), \\
\tau(\alpha) \psi(w)=\psi(\alpha w) & \left(\alpha \in G U\left(\eta_{n}\right), 0 \ll \nu(\alpha) \in F\right), \\
\mu_{v}(\tau(\alpha), \psi(w))=\widehat{A}_{v} \operatorname{diag}\left[\lambda_{v}(\alpha, w), \mu_{v}(\alpha, w)\right] A_{v}^{*} . \tag{6.22c}
\end{array}
$$

Taking a Z-lattice $L_{K}$ in $K_{2 n}^{1}$ such that $\operatorname{Tr}_{K / \mathbf{Q}}\left(x \eta_{n} y^{*}\right) \in \mathbf{Z}$ for every $x, y \in L_{K}$, we consider a family of polarized abelian varieties in Case UT defined in §4.7:

$$
\begin{gather*}
\mathcal{F}\left(\Omega_{K}\right)=\left\{\mathcal{P}_{w} \mid w \in \mathcal{H}_{n}^{\text {a }}\right\}, \quad \mathcal{P}_{w}=\left(A_{w}, \mathcal{C}_{w}, \iota_{w} ;\left\{t_{i}(w)\right\}_{i=1}^{s}\right),  \tag{6.23a}\\
\Omega_{K}=\left\{K, \Psi, L_{K}, \eta_{n},\left\{h\left(u_{i}\right)\right\}_{i=1}^{s}\right\} \tag{6.23b}
\end{gather*}
$$

Take the map $p_{z}$ of (4.24) for $z \in \mathfrak{H}_{2 n}^{\mathbf{a}}$ and similarly $p_{w}$ for $w \in \mathcal{H}_{n}^{\mathbf{a}}$. Then a simple calculation shows that the upper half of $X_{v}\left(p_{w} \circ h\right)$ is of the form

$$
\left[\begin{array}{cccc}
w & \zeta w & 2^{-1} 1_{n} & -(2 \zeta)^{-1} 1_{n}  \tag{6.24a}\\
{ }^{t} w & -\zeta \cdot{ }^{t} w & 2^{-1} 1_{n} & (2 \zeta)^{-1} 1_{n}
\end{array}\right]_{v}
$$

This equals $\bar{A}_{v}^{-1}\left[\begin{array}{ll}\psi(w)_{v} & 1_{2 n}\end{array}\right]$. We also see that

$$
\begin{equation*}
\bar{A} \cdot p_{w}(h(a))=p_{\psi(w)}(a) \quad\left(a \in F_{4 n}^{1}\right) \tag{6.24b}
\end{equation*}
$$

where $\bar{A}=\operatorname{diag}\left[\bar{A}_{v}\right]_{v \in \mathbf{a}}$. Define a Z-lattice $L_{F}$ in $F_{4 n}^{1}$ by $L_{F}=h^{-1}\left(L_{K}\right)$. Then the map $u \mapsto \bar{A} u$ for $u \in\left(\mathbf{C}^{2 n}\right)^{\mathbf{a}}$ gives an isomorphism of $\left(\mathbf{C}^{2 n}\right)^{\mathbf{a}} / p_{w}\left(L_{K}\right)$ onto $\left(\mathbf{C}^{2 n}\right)^{\mathbf{a}} / p_{\psi(w)}\left(L_{F}\right)$, or rather, an isomorphism of $\mathcal{P}_{w}$ to $\mathcal{P}_{\psi(w)}$ if we restrict $\iota_{w}$ to $F$.
6.6. In Case UT we identify $\left(\mathbf{C}^{n}\right)^{\mathbf{b}}$ with $\left(\mathbf{C}^{n}\right)^{\mathbf{a}} \times\left(\mathbf{C}^{n}\right)^{\mathbf{a}}$ via the map $\left(x_{v}\right)_{v \in \mathbf{b}} \mapsto$ $\left(x_{v \rho}\right)_{v \in \mathbf{a}} \times\left(x_{v}\right)_{v \in \mathbf{a}}$. We now define theta functions $\theta_{K}$ and $\varphi_{K}$ by

$$
\begin{align*}
& \theta_{K}(u, w ; \ell)=\sum_{g \in K^{n}} \ell(g) \mathbf{e}_{\mathbf{a}}\left({ }^{t} g w \bar{g}+{ }^{t} g x+g^{*} y\right),  \tag{6.25}\\
& \varphi_{K}(u, w ; \ell)=\mathbf{e}_{\mathbf{a}}\left({ }^{t} y\left(w-w^{*}\right)^{-1} x\right) \theta_{K}(u ; w ; \ell),  \tag{6.26}\\
& \\
& \qquad\left(u=\left[\begin{array}{l}
x \\
y
\end{array}\right] \in\left(\mathbf{C}^{n}\right)^{\mathbf{b}}, w \in \mathcal{H}_{n}^{\mathbf{a}}, \ell \in \mathcal{S}\left(K_{\mathbf{h}}^{n}\right)\right) .
\end{align*}
$$

By an easy calculation we can verify that

$$
\begin{equation*}
\varphi_{K}(u, w ; \ell)=\varphi_{F}(\bar{A} u, \psi(w) ; \ell \circ p) \quad\left(\ell \in \mathcal{S}\left(K_{\mathbf{h}}^{n}\right)\right) \tag{6.27}
\end{equation*}
$$

where $p: F^{2 n} \rightarrow K^{n}$ is defined by $p\left(\left[\begin{array}{l}a \\ b\end{array}\right]\right)=a-\zeta b$ for $a, b \in F^{2 n}$, and the same type of equality holds with $\theta$ in place of $\varphi$.
6.7. In Case SP we have $G=S p(n, F)$. We define, in this case, open subgroups $C^{0}$ and $C^{\theta}$ of $G_{\mathbf{A}}$ and arithmetic subgroups $\Gamma^{0}$ and $\Gamma^{\theta}$ of $G$ by

$$
\begin{gather*}
C^{0}=\left\{\xi \in G_{\mathbf{A}} \mid a_{\xi} \prec \mathfrak{g}, b_{\xi} \prec \mathfrak{d}^{-1}, c_{\xi} \prec \mathfrak{d}, d_{\xi} \prec \mathfrak{g}\right\},  \tag{6.28}\\
C^{\theta}=\left\{\xi \in C^{0} \mid\left(a_{\xi} \cdot{ }^{t} b_{\xi}\right)_{i i} \prec 2 \mathfrak{d}^{-1}\right. \text { and }  \tag{6.29}\\
\left.\quad\left(c_{\xi} \cdot{ }^{t} d_{\xi}\right)_{i i} \prec 2 \mathfrak{d} \text { for } 1 \leq i \leq n\right\}, \\
\Gamma^{0}=G \cap C^{0}, \quad \Gamma^{\theta}=G \cap C^{\theta}, \tag{6.30}
\end{gather*}
$$

where $\mathfrak{d}$ is the different of $F$ relative to $\mathbf{Q}$. Notice that $C^{0}=G_{\mathbf{A}} \cap C\left[\mathfrak{d}^{-1}, \mathfrak{d}\right]$ with $C[$,$] defined by (1.17) with m=n$, and so $C^{0}$ is a subgroup of $G_{\mathbf{A}}$. That $C^{\theta}$ is indeed a subgroup will be shown in §A2.3. We are going to define a factor of automorphy of weight $1 / 2$ in Case SP. To make our formulas short, for $\alpha \in \widetilde{G}_{+}$ and $\kappa \in \mathbf{Z}$ we put

$$
\begin{equation*}
j_{\alpha}(z)^{\kappa \mathbf{a}}=\prod_{v \in \mathbf{a}} j_{v}(\alpha, z)^{\kappa} \quad(\text { Cases SP and UT }) \tag{6.31}
\end{equation*}
$$

with $j_{v}$ of (5.3), in accordance with (5.4b). It should be noted that we use a, not $\mathbf{b}$; naturally we write this $j_{\alpha}(z)^{\mathbf{a}}$ if $\kappa=1$.
6.8. Theorem (Case SP). There is a holomorphic function $h_{\alpha}(z)$, written also $h(\alpha, z)$, on $\mathfrak{H}_{n}^{\mathbf{a}}$ defined for each $\alpha \in \Gamma^{\theta}$ with the following properties:
(1) $h_{\alpha}(z)^{2}=\zeta_{\alpha} j_{\alpha}(z)^{\mathbf{a}}$ with a root of unity $\zeta_{\alpha}$.
(2) $h_{\alpha \beta}(z)=h_{\alpha}(\beta z) h_{\beta}(z)$ for every $\alpha$ and $\beta$ in $\Gamma^{\theta}$.
(3) $h_{\gamma}(z)=1$ if $c_{\gamma}=0$.
(4) Given $\lambda \in \mathcal{S}\left(F_{\mathrm{h}}^{n}\right)$, there is an open subgroup $D_{\lambda}$ of $C^{\theta}$ such that

$$
\begin{equation*}
\varphi_{F}\left({ }^{t} \mu_{\alpha}(z)^{-1} u, \alpha z ; \lambda\right)=h_{\alpha}(z) \varphi_{F}\left(u, z ; \lambda^{\prime}\right) \tag{6.32}
\end{equation*}
$$

with $\lambda^{\prime}(x)=\lambda(d x)$ if $\alpha \in G \cap\left(D_{\lambda} \operatorname{diag}\left[{ }^{t} d^{-1}, d\right]\right)$ with $d \in \prod_{v \in \mathbf{h}} G L_{n}\left(\mathfrak{g}_{v}\right)$, where $\mu_{\alpha}$ is defined by (5.1). In particular, we can take $D_{\lambda}=C^{\theta}$ if $\lambda$ is the characteristic function of $\mathfrak{g}^{n}$.
(5) Let $\varepsilon$ be the Hecke character of $F$ corresponding to the quadratic extension $F(\sqrt{-1}) / F$ and let $\mathfrak{d}$ be the different of $F$ relative to $\mathbf{Q}$. If $\gamma \in \Gamma^{0}, b_{\gamma} \prec 2 \mathfrak{d}^{-1}$, and $c_{\gamma} \prec 2 \mathfrak{2}$, then $\operatorname{det}\left(d_{\gamma}\right)$ is prime to 2 , and

$$
h_{\gamma}(z)^{2}=\prod_{v \mid 2} \varepsilon_{v}\left(\operatorname{det}\left(d_{\gamma}\right)\right) j_{\gamma}(z)^{\mathbf{a}}
$$

This theorem, as well as the following one, will be proven in §A2.9. The function $h_{\gamma}$ may be called the factor of automorphy of weight $\mathbf{a} / 2$ (or simply, of weight $1 / 2$ ).

For the moment, it is defined only for $\gamma \in \Gamma^{\theta}$. We cannot define such for every $\gamma \in S p(n, F)$ consistently so that property (2) above holds in general. Howevere, we can define $h_{\gamma}$ for $\gamma$ in a certain set larger than $\Gamma^{\theta}$, as will be shown in Theorem A2.4 in the Appendix.
6.9. Theorem (Case SP). (1) Given $\alpha \in \widetilde{G}_{+}$, let $r(z)$ be a holomorphic function on $\mathfrak{H}_{n}^{\mathbf{a}}$ such that $r(z)^{2}=\zeta \cdot j_{\alpha}(z)^{\mathbf{a}}$ with a constant $\zeta \in \mathbf{C}^{\times}$. Then there is a congruence subgroup $\Delta$ of $\Gamma^{\theta}$ depending on $(\alpha, r)$ such that $\alpha \Delta \alpha^{-1} \subset \Gamma^{\theta}$ and

$$
h\left(\alpha \gamma \alpha^{-1}, \alpha z\right) r(z)=r(\gamma z) h_{\gamma}(z) \quad \text { for every } \quad \gamma \in \Delta
$$

(2) Given $\alpha$ and $r(z)$ as in (1), suppose $\alpha \in G$. Then for every $\lambda \in \mathcal{S}\left(F_{\mathrm{h}}^{n}\right)$ we have

$$
\varphi_{F}\left({ }^{t} \mu_{\alpha}(z)^{-1} u, \alpha z ; \lambda\right)=r(z) \varphi_{F}\left(u, z ; \lambda^{\prime}\right)
$$

with an element $\lambda^{\prime} \in \mathcal{S}\left(F_{\mathrm{h}}^{n}\right)$, which is determined by $\alpha, r$, and $\lambda$. In particular, if $\alpha=\eta_{n}$ and $r(z)=\prod_{v \in \mathbf{a}} \operatorname{det}\left(-i z_{v}\right)^{1 / 2}$, then $\lambda^{\prime}$ is given by

$$
\begin{equation*}
\lambda^{\prime}(x)=\left|D_{F}\right|^{-n / 2} \int_{\mathbf{F}_{\mathbf{h}}^{n}} \lambda(y) \mathbf{e}_{\mathbf{h}}\left({ }^{t} x y\right) d y . \tag{6.33}
\end{equation*}
$$

Here the branch of $r$ is chosen so that $r(z)>0$ if $\operatorname{Re}(z)=0, D_{F}$ is the discriminant of $F, \mathbf{e}_{\mathbf{h}}$ is the character of $F_{\mathbf{h}}$ defined in $\S 1.6$, and $d y$ is the Haar measure such that the measure of $\prod_{v \in \mathbf{h}} \mathfrak{g}_{v}^{n}$ is 1 .

Put $\Gamma_{\lambda}=G \cap D_{\lambda}$ with $D_{\lambda}$ of Theorem 6.8 (4). Taking $d=1_{n}$ in Theorem 6.8 (4), we obtain

$$
\begin{equation*}
\varphi_{F}\left({ }^{t} \mu_{\gamma}(z)^{-1} u, \gamma z ; \lambda\right)=h_{\gamma}(z) \varphi_{F}(u, z ; \lambda) \quad \text { for every } \quad \gamma \in \Gamma_{\lambda} . \tag{6.34}
\end{equation*}
$$

Now let $s \in F, \gg 0$; take $\alpha=\operatorname{diag}\left[s 1_{n}, 1_{n}\right]$ and $r(z)=1$ in Theorem 6.9 (1); put $\beta=\alpha \gamma \alpha^{-1}$. Then $\mu_{\beta}(s z)=\mu_{\gamma}(z)$ and $h_{\beta}(s z)=h_{\gamma}(z)$ for $\gamma$ in a suitable $\Delta$. Therefore from (6.34) with $\beta$ in place of $\gamma$, we obtain

$$
\begin{equation*}
\varphi_{F}\left({ }^{t} \mu_{\gamma}(z)^{-1} u, s \cdot(\gamma z) ; \lambda\right)=h_{\gamma}(z) \varphi_{F}(u, s z ; \lambda) \quad \text { for every } \quad \gamma \in \Gamma \tag{6.35}
\end{equation*}
$$

with a congruence subgroup $\Gamma$ of $G$ depending on $\lambda$ and $s$.
6.10. Let us now define, in Case SP, modular forms of half-integral weight. First, by an integral weight we mean an element of $\mathbf{Z}^{\mathbf{a}}$. By a half-integral weight we mean an element $k=\left(k_{v}\right)_{v \in \mathbf{a}}$ of $2^{-1} \mathbf{Z}^{\mathbf{a}}$ such that $k_{v}=m_{v}+(1 / 2)$ with $m_{v} \in \mathbf{Z}$ for every $v \in \mathbf{a}$. Given such a $k, \alpha \in \Gamma^{\theta}$, and a $\mathbf{C}$-valued function $f$ on $\mathfrak{H}_{n}^{\mathbf{a}}$, we define a $\mathbf{C}$-valued function $f \|_{k} \alpha$ on $\mathfrak{H}_{n}^{\text {a }}$ by

$$
\begin{equation*}
\left(f \|_{k} \alpha\right)(z)=h_{\alpha}(z)^{-1} j_{\alpha}(z)^{-m} f(\alpha z) . \tag{6.36}
\end{equation*}
$$

For a congruence subgroup $\Gamma$ of $G$ contained in $\Gamma^{\theta}$ we denote by $\mathcal{M}_{k}(\Gamma)$ the set of all holomorphic functions $f$ on $\mathfrak{H}_{n}^{\text {a }}$ which satisfy $f \|_{k} \gamma=f$ for every $\gamma \in \Gamma$ (and also the cusp condition if $n=1$ and $F=\mathbf{Q}$ ). We then denote by $\mathcal{M}_{k}$ the union of $\mathcal{M}_{k}(\Gamma)$ for all such $\Gamma$, and define $\mathcal{A}_{k}$ and $\mathcal{A}_{k}(\Gamma)$ in the same manner as in §5.3. If $f \in \mathcal{M}_{k}$, by Theorem 6.8 (3) we have $f(z+\sigma)=f(z)$ for every $\sigma$ in a suitable lattice in $S$ and $f\left(a x \cdot{ }^{t} a\right)=\operatorname{det}(a)^{-m} f(z)$ for every $a$ in a suitable subgroup of $G L_{n}(\mathfrak{g})$ of finite index. Thus Proposition 5.7 is applicable to $f$, and hence $f$ has an expansion of type (5.27). Now, given a subfield $D$ of $\mathbf{C}$, we can define $\mathcal{M}_{k}(D), \mathcal{A}_{k}(D), \mathcal{M}_{k}(\Gamma, D)$, and $\mathcal{A}_{k}(\Gamma, D)$ in the same manner as in $\S 5.8$. Notice that (5.30) is true for the present $\mathcal{M}_{k}(D)$ too.

Suppose $k \in \mathbf{Z}^{\mathbf{a}}$ and $k_{v}=\kappa$ for every $v \in \mathbf{a}$ with $\kappa \in 2^{-1} \mathbf{Z}$ in Case SP and $\kappa \in \mathbf{Z}$ in Case UT. (This means that $j_{\alpha}(z)^{k}=j_{\alpha}(z)^{\kappa \mathbf{a}}$ if $\kappa \in \mathbf{Z}$; also $k_{v}=0$ for $v \in \mathbf{b}, \notin \mathbf{a}$ in Case UT.) We then employ $\kappa \mathbf{a}$ instead of the subscript $k$ (that is, we write $f \|_{\kappa \mathbf{a}} \alpha$ and $\mathcal{M}_{\kappa \mathbf{a}}(D)$ for $f \|_{k} \alpha$ and $\mathcal{M}_{k}(D)$, for example).

In Case UT we can formulate similar results in a more clear-cut way. Indeed we have:
6.11. Theorem (Case UT). Every element $\alpha$ of $G$ gives a C-linear automorphism of $\mathcal{S}\left(K_{\mathrm{h}}^{n}\right)$, written $\ell \mapsto \ell^{\alpha}$ for $\ell \in \mathcal{S}\left(K_{\mathrm{h}}^{n}\right)$, with the following properties:
(1) $\varphi_{K}\left({ }^{t} M_{\alpha}(z)^{-1} u, \alpha z ; \ell\right)=j_{\alpha}(z)^{\mathbf{a}} \varphi_{K}\left(u, z ; \ell^{\alpha}\right)$, where $M_{\alpha}$ is defined by (5.1) and $j_{\alpha}(z)^{\mathrm{a}}$ by (6.31).
(2) $\ell^{\alpha \beta}=\left(\ell^{\alpha}\right)^{\beta}$.
(3) For every $\ell \in \mathcal{S}\left(K_{\mathbf{h}}^{n}\right)$ there exists a congruence subgroup $\Gamma_{\ell}$ of $G$ such that $\ell^{\gamma}=\ell$ for every $\gamma \in \Gamma_{\ell}$.
(4) In particular, we have

$$
\left(\ell^{\eta}\right)(x)=(-i)^{n[F: \mathbf{Q}]}\left|D_{K}\right|^{-n / 2} \int_{\mathbf{K}_{\mathbf{h}}^{n}} \ell(y) \mathbf{e}_{\mathbf{h}}\left(\operatorname{Tr}_{K / F}\left(y^{*} x\right)\right) d y
$$

where $D_{K}$ is the discriminant of $K$, $\mathbf{e}_{\mathbf{h}}$ is the character of $F_{\mathbf{h}}$ defined in $\S 1.6$, and $d y$ is the Haar measure such that the measure of $\prod_{v \in \mathbf{h}} \mathfrak{r}_{v}^{n}$ is 1 .

Proof. This is essentially a special case of [S97, Theorem A7.4]. Indeed, define $f\left(z ; u, u^{\prime} ; \ell\right)$ of $[\mathrm{S} 97,(\mathrm{~A} 7.3 .2)]$ with $q=1$ and $H=Q=A=1$. Then we can easily verify that

$$
\varphi_{K}\left(\left[\begin{array}{l}
x  \tag{6.37}\\
y
\end{array}\right], z ; \ell\right)=f(2 z ; y, x ; \ell)
$$

In [S97, Theorem A7.4] we proved the results of type (1,2,3) of our theorem for the function $f$ in a stronger form for the elements of the group $G_{1}=G \cap S L_{2 n}(K)$. Given $\alpha \in G_{1}$, put $\beta=\xi^{-1} \alpha^{-1} \xi$ with $\xi=\operatorname{diag}\left[1_{n}, 2 \cdot 1_{n}\right]$; write $\ell^{\alpha}$ for the symbol ${ }^{\beta} \ell$ defined in that theorem. Then we obtain the present theorem at least for the elements of $G_{1}$. Now $G$ is generated by $G_{1}$ and the elements of the form $\operatorname{diag}[a, \widehat{a}]$ with $a \in G L_{n}(K)$. Therefore we can easily extend the results to $G$. As for (4), if $\alpha=\eta$, then $\beta=\tau \eta$ with $\tau=\operatorname{diag}\left[-2,-2^{-1}\right]$, so that $\ell^{\eta}={ }^{\beta} \ell={ }^{\tau}\left({ }^{\eta} \ell\right)$. Combining (5) and (6) of [S97, Theorem A7.4], we obtain $\ell^{\eta}$ as stated in (4).

Notice that in Case UT the facts corresponding to Theorem 6.9 (1) and (6.35) (with $j_{\gamma}^{\text {a }}$ in place of $h_{\gamma}$ ) are trivial. Also, we shall later give Theorem A5.4 which essentially includes Theorem 6.11, as well as [S97, Theorem A7.4].
6.12. Theorem. Let $\mathcal{F}\left(\Omega_{K}\right)$ be defined in Cases $S P$ and $U T$ as in $\S \S 6.3$ and 6.5. Then there exist a finite set $\Lambda$ of $\mathbf{Z}$-valued elements of $\mathcal{S}\left(K_{\mathrm{h}}^{n}\right)$, a positive integer $p$, and a congruence subgroup $\Gamma_{0}$ of $G$ with the following properties, in which it is understood that $K=F$ in Case SP:
(1) The quotient $\theta_{K}\left(u, p^{-1} z ; \lambda\right) / \theta_{K}\left(u, p^{-1} z ; \lambda^{\prime}\right)$ is invariant under $u \mapsto u+\ell$ for every $\lambda, \lambda^{\prime} \in \Lambda$ and every $\ell \in p_{z}\left(L_{K}\right)$.
(2) For every $\left(u_{0}, z_{0}\right) \in\left(\mathbf{C}^{n}\right)^{\mathbf{b}} \times \mathcal{H}$ there exists an element $\lambda$ of $\Lambda$ such that $\theta_{K}\left(u_{0}, p^{-1} z_{0} ; \lambda\right) \neq 0$. Let $\Theta_{K}(u ; z)$ denote the point in the complex projective space $P^{m}(\mathbf{C})$ whose homogeneous coordinates are $\left(\theta_{K}\left(u, p^{-1} z ; \lambda\right)\right)_{\lambda \in \Lambda}$, where $m=$ $\#(\Lambda)-1$.
(3) For each fixed $z \in \mathcal{H}$ the map $u \mapsto \Theta_{K}(u, z)$ defines a biregular projective embedding of $\left(\mathbf{C}^{n}\right)^{\mathbf{b}} / p_{z}\left(L_{K}\right)$ onto an abelian variety.
(4) The polarization on the image abelian variety determined by its hyperplane sections corresponds to the Riemann form of (4.13).
(5) $\Theta_{K}\left({ }^{t} M_{\gamma}(z) u, z\right)=\Theta_{K}(u, \gamma z)$ for every $\gamma \in \Gamma_{0}$.

Proof. We first consider $\mathbf{C}^{n} / p_{z}(L)$ for $z \in \mathfrak{H}_{n}$ and a Z-lattice $L$ in $\mathbf{Q}_{2 n}^{1}$ such that $x \eta_{n} \cdot{ }^{t} y \in \mathbf{Z}$ for every $x, y \in L$. As is well-known, we can find an element $\alpha$ of $\mathbf{Z}_{2 n}^{2 n}$ such that $\mathbf{Z}_{2 n}^{1} \alpha=L$ and $\alpha \eta_{n} \cdot{ }^{t} \alpha=\left[\begin{array}{cc}0 & -\delta \\ \delta & 0\end{array}\right]$ with a diagonal matrix $\delta$. Put $\tau_{\delta}=\operatorname{diag}\left[1_{n}, \delta\right]$. Then $\alpha^{-1} \tau_{\delta}$ gives an isomorphism of $\left(\eta_{n}, L\right)$ onto $\left(\eta_{n}, \mathbf{Z}_{2 n}^{1} \tau_{\delta}\right)$. Thus we may assume, without losing generality, that $L=\mathbf{Z}_{2 n}^{1} \tau_{\delta}$. Then $p_{z}(L)=$ $\left\{z a+\delta b \mid a, b \in \mathbf{Z}^{n}\right\}$. Take an integer $p$ so that every entry of $p \delta$ is divisible by an integer greater than 2 , and put $f_{r}(u)=\varphi(p u, p z ; r, 0)$ for $r \in(p \delta)^{-1} \mathbf{Z}^{n}$. Let $\mathfrak{T}$ be the vector space of all the holomorphic functions $f$ on $\mathbf{C}^{n}$ such that

$$
\begin{aligned}
& f(u+\ell)=\mathbf{e}\left((p / 2) \cdot{ }^{t} a \delta b\right) \mathbf{e}\left(p \ell^{*}(z-\bar{z})^{-1}(u+(\ell / 2)) f(u)\right. \\
& \quad \text { for every } \ell=z a+\delta b \text { with } a, b \in \mathbf{Z}^{n} .
\end{aligned}
$$

(This is essentially the same as (2.4).) Let $R$ be a complete set of representatives for $(p \delta)^{-1} \mathbf{Z}^{n} / \mathbf{Z}^{n}$. Then $\left\{f_{r} \mid r \in R\right\}$ is a C-basis of $\mathfrak{T}$, and the map $u \mapsto\left(f_{r}(u)\right)_{r \in R}$ defines a biregular projective embedding of $\mathbf{C}^{n} / p_{z}(L)$. These are well-known classical results. (For detailed treatments, see [W58] and [S98, Section 27]. Observe that $u \mapsto p u$ gives an isomorphism of $\mathbf{C}^{n} / p_{z}\left(\mathbf{Z}_{2 n}^{1} \tau_{\delta}\right)$ onto $\mathbf{C}^{n} / p_{p z}\left(\mathbf{Z}_{2 n}^{1} \tau_{p \delta}\right)$; apply the standard facts as stated in $[\mathrm{S} 98, \S 27.12]$ to the last complex torus.) Since the Riemann form corresponds to the zero divisor of a nonzero $f$, it corresponds to the hyperplane sections of the image variety. Now we obtain the desired $\Lambda$ as follows. In Case SP, we may assume that $g^{-1}\left(L_{F}\right)=\mathbf{Z}_{2 e n}^{1} \tau_{\delta}$ with $\delta$ of size en. Then we consider $\varphi_{F}(p u, p z ; r, 0)$ with $r \in g_{1}\left((p \delta)^{-1} \mathbf{Z}^{e n}\right) / g_{1}\left(\mathbf{Z}^{e n}\right)$. We can express this as $\varphi_{F}\left(u, p^{-1} z ; \lambda_{r}\right)$ with some $\lambda_{r} \in \mathcal{S}\left(F_{\mathbf{h}}^{n}\right)$ that is $\mathbf{Z}$-valued. We then take $\Lambda=\left\{\lambda_{r}\right\}$ with all such $r$ 's. Since $u \mapsto B u$ for $u \in\left(\mathbf{C}^{n}\right)^{\mathbf{a}}$ gives an isomorphism of $\left(\mathbf{C}^{n}\right)^{\mathbf{a}} / p_{z}(L)$ onto $\mathbf{C}^{e n} / p_{\varepsilon(z)}\left(L_{0}\right)$, from (6.15) and the above classical results we obtain the first four assertions in Case SP. The existence of $\Gamma_{0}$ satisfying (5) follows from (6.35). Case UT can be handled in a similar manner by means of (6.27) and what we said at the end of $\S 6.5$.
6.13. With a fixed choice of $\Lambda$ let $A_{w}$ denote the image abelian variety of the map of Theorem 6.12(3). Hereafter we understand that the symbol $A_{w}$ in $\mathcal{P}_{w}$ belonging to our families $\mathcal{F}\left(\Omega_{F}\right)$ and $\mathcal{F}\left(\Omega_{K}\right)$ is this projective variety. From Theorem 6.12 (5) we see that

is a commutative diagram, where $\Theta_{z}(u)=\Theta_{K}(u, z)$, and $K=F$ in Case SP.
6.14. Proposition. Case SP: Let $\ell \in \mathcal{S}\left(F_{\mathrm{h}}^{n}\right), b, c \in F^{n}$, and $r \in F, \gg 0$. Define $\ell^{\prime} \in \mathcal{S}\left(F_{\mathbf{h}}^{n}\right)$ by $\ell^{\prime}(h)=\ell(h-b) \mathbf{e}_{\mathbf{a}}\left({ }^{t} c(h-b)\right)$. Then

$$
\begin{equation*}
\mathbf{e}_{\mathbf{a}}\left((r / 2) \cdot{ }^{t} b z b\right) \theta_{F}(r z b+c, r z ; \ell)=\theta_{F}\left(0, r z ; \ell^{\prime}\right) \tag{6.39}
\end{equation*}
$$

Moreover, $\theta_{F}(0, r z ; \ell)$ as a function of $z$ belongs to $\mathcal{M}_{\mathbf{a} / 2} ;$ it belongs to $\mathcal{M}_{\mathbf{a} / 2}(D)$ if $\ell$ is $D$-valued for any subfield $D$ of $\mathbf{C}$.

Case UT: Let $\ell \in \mathcal{S}\left(K_{\mathrm{h}}^{n}\right), b, c \in K^{n}$, and $r \in F, \gg 0$. Define $\ell^{\prime} \in \mathcal{S}\left(K_{\mathrm{h}}^{n}\right)$ by $\ell^{\prime}(h)=\ell(h-b) \mathbf{e}\left(\operatorname{Tr}_{K / \mathbf{Q}}\left(c^{*}(h-b)\right)\right)$. Then

$$
\begin{equation*}
\mathbf{e}_{\mathbf{a}}\left(r \cdot{ }^{t} b w \bar{b}\right) \theta_{K}\left(r w \bar{b}+\bar{c}, r \cdot{ }^{t} w b+c ; r w ; \ell\right)=\theta_{K}\left(0, r w ; \ell^{\prime}\right) \tag{6.40}
\end{equation*}
$$

Moreover, $\theta_{K}(0, r z ; \ell)$ as a function of $z$ belongs to $\mathcal{M}_{\mathbf{a}}$; it belongs to $\mathcal{M}_{\mathbf{a}}(D)$ if $\ell$ is $D$-valued for any subfield $D$ of $\mathbf{C}$.

Proof. Since $\theta_{F}(0 ; z ; \ell)=\varphi_{F}(0 ; z ; \ell)$, we see from (6.35) that $\theta_{F}(0, r z ; \ell) \in$ $\mathcal{M}_{\mathrm{a} / 2}$ for any $\ell$. The $D$-rationality for $D$-valued $\ell$ is obvious. Equality (6.39) can be verified by a direct calculation. Case UT is similar; we need Theorem 6.11 ( 1 , $3)$ instead of (6.35).
6.15. Proposition (Case SP). Let $f_{i} \in \mathcal{A}_{k_{i}}$ with an integral or a half-integral weight $k_{i}$ for $1 \leq i \leq m$. Then $f_{1} \cdots f_{m} \in \mathcal{A}_{k}$ with $k=\sum_{i=1}^{m} k_{i}$.

Proof. From Theorem 6.8 (5) we see that $h_{\gamma}(z)^{2}=j_{\gamma}(z)^{\mathbf{a}}$ if $\gamma \in \Gamma^{0}, b_{\gamma} \prec$ $2 \mathfrak{d}^{-1}, c_{\gamma} \prec 2 \mathfrak{d}$, and $\operatorname{det}\left(d_{\gamma}\right)-1 \in 4 \mathfrak{g}$. Our assertion can easily be derived from this fact.
6.16. Proposition. For an integral or a half-integral weight $k$ let $0 \neq f(z)=$ $\sum_{h \in S} c(h) \mathbf{e}_{\mathbf{a}}^{n}(h z) \in \mathcal{M}_{k}$ and let $r=\operatorname{Max}\{\operatorname{rank}(h) \mid c(h) \neq 0\}$. Then

$$
\begin{align*}
& r=n \Longleftrightarrow\left\{\begin{array}{lll}
k_{v} \geq n / 2 & \text { for every } v \in \mathbf{a} & \text { (Case SP), } \\
k_{v \rho}+k_{v} \geq n & \text { for every } v \in \mathbf{a} & \text { (Case UT); }
\end{array}\right.  \tag{6.41a}\\
& r<n \Longleftrightarrow\left\{\begin{array}{lll}
k_{v}=r / 2<n / 2 & \text { for every } v \in \mathbf{a} & \text { (Case SP), } \\
k_{v \rho}+k_{v}=r<n & \text { for every } v \in \mathbf{a} & \text { (Case UT). }
\end{array}\right.
\end{align*}
$$

For the proof, see [S94b, Theorem 5.6 and Corollary 5.7]. Clearly it follows that $\mathcal{M}_{0}=\mathbf{C}$ in both cases. If $0 \neq f \in \mathcal{S}_{k}$, then $r=n$, and hence

$$
\mathcal{S}_{k} \neq\{0\} \Longrightarrow\left\{\begin{array}{ll}
k_{v} \geq n / 2 & \text { for every } v \in \mathbf{a}  \tag{6.42}\\
k_{v \rho}+k_{v} \geq n & \text { for every } v \in \mathbf{a}
\end{array} \quad\right. \text { (Case SP) }
$$

6.17. Lemma. Let $0<\kappa \in 2^{-1} \mathbf{Z}$ in Case SP and let $0<\kappa \in \mathbf{Z}$ in Case UT. Then, given $z_{0} \in \mathcal{H}$, there exists an element $f \in \mathcal{M}_{\kappa \mathbf{a}}(\mathbf{Q})$ such that $f\left(z_{0}\right) \neq 0$.

Proof. By Theorem 6.12 (2) there is a $\mathbf{Z}$-valued $\lambda$ such that $\theta_{K}\left(0, p^{-1} z_{0} ; \lambda\right)$ $\neq 0$. Put $g(z)=\theta_{K}\left(0, p^{-1} z ; \lambda\right)$. By Proposition 6.14, $g \in \mathcal{M}_{\mathbf{a} / 2}(\mathbf{Q})$ in Case SP and $g \in \mathcal{M}_{\mathbf{a}}(\mathbf{Q})$ in Case UT. Thus a suitable power of $g$ gives the desired $f$.

## CHAPTER II

## ARITHMETICITY OF AUTOMORPHIC FORMS

## 7. The field $\mathcal{A}_{0}\left(\mathrm{Q}_{\mathrm{ab}}\right)$

The principal result of this section is Theorem 7.10. We start with some auxiliarly lemmas.
7.1. Lemma. Let $h$ and $k$ be subfields of $\mathbf{C}$ with countably many elements. Then the following assertions hold:
(i) If $k$ is stable under $\operatorname{Aut}(\mathbf{C} / h)$, then the composite field $h k$ is a finite or an infinite Galois extension of $h$.
(ii) If every element of $\operatorname{Aut}(\mathbf{C} / h)$ gives the identity map on $k$, then $k \subset h$.
(iii) If $Y$ is an algebro-geometric object such as a variety, a divisor, or a rational map, and $Y^{\sigma}=Y$ for every $\sigma \in \operatorname{Aut}(\mathbf{C} / h)$, then $Y$ is rational over $h$.

The proof, being completely elementary, is left to the reader. In the following sections we shall often make use of these principles, though we shall not explicitly mention them in each instance.
7.2. Lemma. Let $\Phi$ and $\Psi$ be extensions of a field $k$ which are linearly disjoint over $k$; let $\Phi^{\prime}$ be a subfield of $\Phi$ containing $k$ such that $\Psi \Phi=\Psi \Phi^{\prime}$. Then $\Phi=\Phi^{\prime}$.

Proof. Let $f \in \Phi$. Then $f \in \Psi \Phi=\Psi \Phi^{\prime}$, so that $f=\sum_{\kappa} a_{\kappa} p_{\kappa} / \sum_{\lambda} b_{\lambda} q_{\lambda}$ with $a_{\kappa}, b_{\lambda} \in \Psi$ and $p_{\kappa}, q_{\lambda} \in \Phi^{\prime}$. Take a finite set $\left\{c_{\mu}\right\}$ of elements of $\Psi$ linearly independent over $k$ so that $\sum_{\kappa} k a_{\kappa}+\sum_{\lambda} k b_{\lambda}=\sum_{\mu} k c_{\mu}$. Expressing $a_{\kappa}$ and $b_{\lambda}$ as $k$-linear combinations of $c_{\mu}$, we can put $f=\sum_{\mu} c_{\mu} g_{\mu} / \sum_{\mu} c_{\mu} h_{\mu}$ with $g_{\mu}, h_{\mu} \in \Phi^{\prime}$. Then $\sum_{\mu} c_{\mu}\left(f h_{\mu}-g_{\mu}\right)=0$. Since $f h_{\mu}-g_{\mu} \in \Phi$, the linear disjointness shows that $f h_{\mu}=g_{\mu}$ for every $\mu$. Thus $f=g_{\mu} / h_{\mu}$ for some $\mu$, and hence $f \in \Phi^{\prime}$ as expected.
7.3. Lemma. Let $\left\{f_{\nu} \mid \nu \in N\right\}$ be a set of meromorphic functions in a connected open subset $D$ of $\mathbf{C}^{d}$, indexed by an at most countable set $N$. Let $k$ be a subfield of $\mathbf{C}$ with only countably many elements. Then there exists a point $z_{0}$ of $D$ such that the specialization $\left\{f_{\nu}\right\}_{\nu \in N} \mapsto\left\{f_{\nu}\left(z_{0}\right)\right\}_{\nu \in N}$ defines an isomorphism of the field $k\left(f_{\nu} \mid \nu \in N\right)$ onto $k\left(f_{\nu}\left(z_{0}\right) \mid \nu \in N\right)$ over $k$.

Proof. We may assume that $N=\{1,2,3, \ldots\}$ (finite or not). By induction we can find a subset $M=\left\{\nu_{1}, \nu_{2}, \ldots\right\}$ of $N$ such that: (i) $\nu_{1}<\nu_{2}<\cdots$; (ii) $f_{\nu_{1}}, f_{\nu_{2}}, \ldots$ are algebraically independent over $k$; and (iii) $f_{1}, \ldots, f_{n}$ are algebraic over $k\left(f_{\nu} \mid \nu \in M, \nu \leq n\right)$. Let $S_{m}$ be the set of all polynomials $P\left(X_{1}, \ldots, X_{m}\right) \neq 0$ in $m$ indeterminates with coefficients in $k$, and $W_{\nu}$ the set of the points of $D$ where $f_{\nu}$ is not holomorphic. For each $P \in S_{m}$ put

$$
E_{P}=\left\{z \in D-\bigcup_{i=1}^{m} W_{\nu_{i}} \mid P\left(f_{\nu_{1}}(z), \ldots, f_{\nu_{m}}(z)\right)=0\right\} .
$$

The closure of $E_{P}$ in $D$ has no interior point of $D$. Now observe that $S_{m}$ has only countably many elements. Recall a well-known fact that if $D$ is covered by countably many closed subsets, then at least one of them has an interior point. Therefore we find a point $z_{0}$ of $D$ not belonging to the countable union $\left[\bigcup_{\nu \in N} W_{\nu}\right] \cup$ $\left[\bigcup_{m=1}^{\infty} \bigcup_{P \in S_{m}} E_{P}\right]$. Then our construction shows that $k\left(f_{1}, \ldots, f_{n}\right)$ has the same transcendence degree as $k\left(f_{1}\left(z_{0}\right), \ldots, f_{n}\left(z_{0}\right)\right)$ over $k$ for every $n$. Therefore the specialization $f_{\nu} \mapsto f_{\nu}\left(z_{0}\right)$ defines an isomorphism of these fields as expected.

We say that $z_{0}$ is generic for $\left\{f_{\nu} \mid \nu \in N\right\}$ over $k$ if $z_{0}$ has the property as in the above lemma. If $z_{0}$ is such a point, $\left\{g_{\lambda} \mid \lambda \in L\right\}$ is another countable set of meromorphic functions on $D$, and each $g_{\lambda}$ is algebraic over the field generated by the $f_{\nu}$ over $k$, then clearly $z_{0}$ is generic for $\left\{f_{\nu} \mid \nu \in N\right\} \cup\left\{g_{\lambda} \mid \lambda \in L\right\}$ over $k$.

Hereafter until the end of Section 8 we treat only Cases SP and UT. We first consider the fields $\mathcal{A}_{0}(k)$ and $\mathcal{A}_{0}(\Gamma, k)$ defined in $\S 5.8$.
7.4. Lemma. (1) For every subfield $k$ of $\mathbf{C}$, the fields $\mathcal{A}_{0}(k)$ and $\mathbf{C}$ are linearly disjoint over $k$.
(2) Let $\Gamma$ be a congruence subgroup of $G$ and let $(V, \varphi)$ be a model of $\Gamma \backslash \mathcal{H}$ in the sense of §4.10. Let $k$ be a subfield of $\mathbf{C}$ such that $V$ is defined over $k$ and $k(V) \circ \varphi \subset \mathcal{A}_{0}(\Gamma, k)$. Then $k(V) \circ \varphi=\mathcal{A}_{0}(\Gamma, k)$, where $k(V)$ is defined as in $\S 2.4$.

Proof. To prove (1), let $f_{1}, \ldots, f_{m}$ be elements of $\mathcal{A}_{0}(k)$ linearly independent over $k$; put $f_{i}=p_{i} / q_{i}$ with $p_{i}, q_{i} \in \mathcal{M}_{\nu_{i}}(k), \nu_{i} \in \mathbf{Z}^{\mathbf{b}}$, and $r_{i}=q_{1} \cdots q_{m} f_{i}$. Then the $r_{i}$ are automorphic forms of the same weight linearly independent over $k$, and hence, linearly independent over $\mathbf{C}$ by virtue of Lemma 5.10 (1). Therefore we obtain (1). As for (2), we have $\mathbf{C} k(V) \circ \varphi \subset \mathbf{C} \mathcal{A}_{0}(\Gamma, k) \subset \mathcal{A}_{0}(\Gamma)=\mathbf{C}(V) \circ \varphi=$ $\mathbf{C} k(V) \circ \varphi$, and hence $\mathbf{C} k(V) \circ \varphi=\mathbf{C} \mathcal{A}_{0}(\Gamma, k)$. Therefore, by (1) and Lemma 7.2 we obtain the desired equality of (2).
7.5. Lemma. Let $P=\left\{\xi \in \widetilde{G} \mid c_{\xi}=0\right\}$, where $c_{\xi}$ is the c-block of $\xi$ (see Lemma 1.9). Put $G_{1}=G \cap S L_{2 n}(K)$ and denote by $\mathcal{G}$ any of the groups $\widetilde{G}, \widetilde{G}_{+}, G$, and $G_{1} ;$ put $\mathcal{P}=P \cap \mathcal{G}$. Then $\mathcal{G}$ is dense in $\mathcal{G}_{\mathbf{a}}$, and $\mathcal{G}$ (resp. $\mathcal{G}_{\mathbf{a}}$ ) is generated by $\eta_{n}$ and $\mathcal{P}\left(\right.$ resp. $\left.\mathcal{P}_{\mathbf{a}}\right)$, where we understand that $\left(\widetilde{G}_{+}\right)_{\mathbf{a}}=\widetilde{G}_{\mathbf{a}+}$ and $\mathcal{P}_{\mathbf{a}}=P_{\mathbf{a}} \cap \mathcal{G}_{\mathbf{a}}$.

Proof. Put $\mathcal{B}=\left\{\xi \in \mathcal{G} \mid \operatorname{det}\left(c_{\xi}\right) \neq 0\right\}$ and $\mathcal{B}_{\mathbf{a}}=\left\{\xi \in \mathcal{G}_{\mathbf{a}} \mid \operatorname{det}\left(c_{\xi}\right) \in K_{\mathbf{a}}^{\times}\right\}$. Then $\mathcal{B}=\mathcal{P} \eta \mathcal{P}$ and $\mathcal{B}_{\mathrm{a}}=\mathcal{P}_{\mathbf{a}} \eta \mathcal{P}_{\mathrm{a}}$. Indeed, we can easily verify that $\mathcal{P} \eta \mathcal{P} \subset \mathcal{B}$ and $\mathcal{P}_{\mathbf{a}} \eta \mathcal{P}_{\mathbf{a}} \subset \mathcal{B}_{\mathbf{a}}$. That $\mathcal{B} \subset \mathcal{P} \eta \mathcal{P}$ and $\mathcal{B}_{\mathbf{a}} \subset \mathcal{P}_{\mathbf{a}} \eta \mathcal{P}_{\mathbf{a}}$ can be seen from an equality

$$
\xi=\left[\begin{array}{ll}
a & b  \tag{*}\\
c & d
\end{array}\right]=\left[\begin{array}{cc}
s \cdot 1_{n} & a c^{-1} \\
0 & 1_{n}
\end{array}\right]\left[\begin{array}{cc}
0 & -1_{n} \\
1_{n} & 0
\end{array}\right]\left[\begin{array}{ll}
c & d \\
0 & \widehat{c}
\end{array}\right]
$$

where $\xi \in G U\left(\eta_{n}\right), s=\nu(\xi)$, and $c$ is invertible. (Notice that by Lemma 1.3 (2), if $\nu(\xi)=\operatorname{det}(\xi)=1$, then $\operatorname{det}(c) \in F$, so that the last matrix of $\left({ }^{*}\right)$ has determinant 1.) Now $\mathcal{B}_{\mathbf{a}}$ is open and dense in $\mathcal{G}_{\mathbf{a}}$. Clearly $\mathcal{P}$ is dense in $\mathcal{P}_{\mathbf{a}}$, and so $\mathcal{B}$ is dense in $\mathcal{B}_{\mathbf{a}}$. Therefore $\mathcal{G}$ is dense in $\mathcal{G}_{\mathbf{a}}$. Now given $\alpha \in \mathcal{G}$, the set $\left\{\xi \in \mathcal{G}_{\mathbf{a}} \mid \operatorname{det}\left(c_{\xi} c_{\alpha \xi}\right) \in K_{\mathbf{a}}^{\times}\right\}$is open in $\mathcal{G}_{\mathbf{a}}$. Therefore we can find $\xi \in \mathcal{G}$ such that $\operatorname{det}\left(c_{\xi} c_{\alpha \xi}\right) \neq 0$. Then both $\xi$ and $\alpha \xi$ belong to $\mathcal{P} \eta \mathcal{P}$, and so $\alpha=\alpha \xi \xi^{-1} \in \mathcal{P} \eta_{n} \mathcal{P} \eta_{n} \mathcal{P}$, which gives the desired fact for $\mathcal{G}$. The assertion for $\mathcal{G}_{\mathrm{a}}$ can be proved in the same manner.
7.6. We now take our setting to be that of Section 6 . We fix a positive integer $p \geq 3$, a subset $\Lambda$ of $\mathcal{S}\left(K_{\mathrm{h}}^{n}\right)$, and a congruence subgroup $\Gamma_{0}$ of $G$ as in Theorem 6.12; we then consider the map $\Theta_{K}$ defined with these $p$ and $\Lambda$ as in (3) of that
theorem. Let $P^{*}(\mathbf{C})$ be the complex projective space in which the map $\Theta_{K}$ takes its values. Then $A_{w}$ of $\S 6.13$ is a subvariety of $P^{*}(\mathbf{C})$. We retain the convention that $K=F$ and $\mathfrak{r}=\mathfrak{g}$ in Case SP (see $\S 3.5$ ). The symbol $p_{z}$ of (4.23) is a map of $\left(K_{\mathrm{a}}\right)_{2 n}^{1}$ into $\left(\mathbf{C}^{n}\right)^{\mathbf{b}}$. Recall the formula $p_{z}(x \alpha)={ }^{t} \mu_{\alpha}(z) p_{\alpha z}(x)$ stated in (4.31). Substituting $p_{\gamma z}(a)$ with $a \in K_{2 n}^{1}$ for $u$ in Theorem 6.12 (5), we obtain

$$
\begin{equation*}
\Theta_{K}\left(p_{z}(a \gamma), z\right)=\Theta_{K}\left(p_{\gamma z}(a), \gamma z\right) \quad \text { for every } a \in K_{2 n}^{1} \text { and } \gamma \in \Gamma_{0} \tag{7.1}
\end{equation*}
$$

Define $t: K_{2 n}^{1} \times \mathcal{H} \rightarrow P^{*}(\mathbf{C})$ and $t_{z}: K_{2 n}^{1} \rightarrow P^{*}(\mathbf{C})$ for $z \in \mathcal{H}$ by

$$
\begin{equation*}
t(a, z)=t_{z}(a)=\Theta_{K}\left(p_{z}(a), z\right)=\left(\theta_{K}\left(p_{z}(a), p^{-1} z ; \lambda\right)\right)_{\lambda \in \Lambda} \tag{7.2}
\end{equation*}
$$

Then from (7.1) we obtain

$$
\begin{equation*}
t(a \gamma, z)=t(a, \gamma z) \quad \text { for every } \quad a \in K_{2 n}^{1} \text { and } \gamma \in \Gamma_{0} \tag{7.3}
\end{equation*}
$$

7.7. We fix an $\mathfrak{r}$-lattice $L$ in $K_{2 n}^{1}$, and for each positive integer $N$ we fix a subset $\left\{u_{i}\right\}_{i=1}^{s}$ of $N^{-1} L$ so that $N^{-1} L=L+\sum_{i=1}^{s} \mathbf{Z} u_{i}$. We then consider the families of $\S \S 6.3$ and 6.5 in the following form:

$$
\begin{gather*}
\mathcal{F}\left(\Omega^{N}\right)=\left\{\mathcal{P}_{z}^{N} \mid z \in \mathcal{H}\right\}, \quad \mathcal{P}^{N}(z)=\mathcal{P}_{z}^{N}=\left(A_{z}, \mathcal{C}_{z}, \iota_{z} ;\left\{t_{i}(z)\right\}_{i=1}^{s}\right)  \tag{7.4a}\\
\Omega^{N}=\left\{K, \Psi, L, \eta_{n},\left\{u_{i}\right\}_{i=1}^{s}\right\} \tag{7.4b}
\end{gather*}
$$

We write $\Omega^{N}$ and $\mathcal{P}_{z}^{N}$ instead of $\Omega$ and $\mathcal{P}_{z}$ in order to emphasize $N$. We simply write $L$ instead of $L_{K}$ in Theorem 6.12. Naturally $t_{i}(z)=t_{z}\left(u_{i}\right)$, and so

$$
\begin{equation*}
\left\{t \in A_{z} \mid N t=0\right\}=\sum_{i=1}^{s} \mathbf{Z} t_{i}(z) \tag{7.4c}
\end{equation*}
$$

For $c \in K$ we have $p_{z}(c a)=\Psi(c) p_{z}(a)$ and $\iota_{z}(c)$ is represented by $\Psi(c)$, and so we obtain

$$
\begin{equation*}
\iota_{z}(c) t_{z}(a)=t_{z}(c a) \quad \text { for every } c \in \mathfrak{r} \text { and } a \in K_{2 n}^{1} \tag{7.5}
\end{equation*}
$$

Define a congruence subgroup $\Gamma^{N}$ of $G$ by

$$
\begin{equation*}
\Gamma^{N}=\{\gamma \in G \mid L \gamma=L, L(\gamma-1) \subset L\} \tag{7.6}
\end{equation*}
$$

By Theorem 4.8, $\mathcal{P}_{z}^{N}$ and $\mathcal{P}_{w}^{N}$ are isomorphic if and only if $z=\gamma w$ for some $\gamma \in \Gamma^{N}$, since $\Gamma$ of (4.28) coincides with $\Gamma^{N}$ for the present family.

Let us now fix a model $\left(V_{N}, \varphi_{N}\right)$ of $\Gamma_{N} \backslash \mathcal{H}$. If $N$ divides $M$, then there exists a rational map $p_{N}^{M}: V_{M} \rightarrow V_{N}$ such that $\varphi_{N}=p_{N}^{M} \circ \varphi_{M}$. We can find a subfield $k_{0}$ of $\mathbf{C}$ with countably many elements over which the varieties $V_{N}$ and the maps $p_{N}^{M}$ are defined for all $M$ and $N$. Put

$$
\begin{equation*}
\mathfrak{F}=\bigcup_{N=1}^{\infty} \mathfrak{F}_{N}, \quad \mathfrak{F}_{N}=\left\{f \circ \varphi_{N} \mid f \in k_{0}\left(V_{N}\right)\right\} \tag{7.7}
\end{equation*}
$$

Then $\mathcal{A}_{0}\left(\Gamma_{N}\right)=\mathbf{C} \mathfrak{F}_{N}$ and $\mathcal{A}_{0}=\mathbf{C} \mathfrak{F}$. Since $\mathfrak{F}_{N}$ is finitely generated over $k_{0}$, we see that $\mathfrak{F}$ is countable. Now let $\mathfrak{K}$ be the field generated over $\mathbf{Q}$ by the functions on $\mathcal{H}$ of the form

$$
\begin{equation*}
\theta_{K}\left(p_{z}(a), p^{-1} z ; \ell\right) / \theta_{K}\left(p_{z}(a), p^{-1} z ; \ell^{\prime}\right) \tag{7.8}
\end{equation*}
$$

for all $\ell, \ell^{\prime} \in \Lambda$ and all $a \in K_{2 n}^{1}$, where $\Lambda$ is the set of Theorem 6.12. By Proposition 6.14 each such quotient belongs to $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right) ;$ thus $\mathfrak{K} \subset \mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$. For each $w \in \mathcal{H}$ let $\mathfrak{K}[w]$ denote the field generated over $\mathbf{Q}$ by the values $f(w)$ for every $f \in \mathfrak{K}$ finite at $w$. Since any affine coordinate of $t_{z}(a)$ is of the form (7.8), $t_{w}(a)$ is rational over $\mathfrak{K}[w]$ for every $a \in K_{2 n}^{1}$, and hence $A_{w}$ is rational over $\mathfrak{K}[w]$. Also, from (7.5)
we see that $\iota_{w}(c)$ is rational over $\mathfrak{K}[w]$. Since the polarization is determined by hyperplane sections, $\mathcal{P}_{w}^{N}$ is rational over $\mathfrak{K}[w]$ for every $N$. Now by Lemma $2.6, \mathbf{Q}_{\mathrm{ab}}$ is contained in the field generated over $\mathbf{Q}$ by the affine coordinates of the points $t_{z}(a)$, and hence $\mathbf{Q}_{\mathrm{ab}} \subset \mathfrak{K}[w]$. Taking $w$ to be generic for the elements of $\mathfrak{K}$ over $\mathbf{Q}$, we see that $\mathbf{Q}_{\mathrm{ab}} \subset \mathfrak{K}$.
7.8. Lemma. Given $w, w^{\prime} \in \mathcal{H}$, suppose that there exist an isomorhism $\sigma$ of $\mathfrak{K}[w]$ onto $\mathfrak{K}\left[w^{\prime}\right]$ and an $\mathfrak{r}$-linear automorphism $\lambda$ of $K_{2 n}^{1} / L$ such that $t(a, w)^{\sigma}=$ $t\left(\lambda(a), w^{\prime}\right)$ for every $a$ and $\lambda\left(u_{i}\right)=u_{i}$ for every $i$. Then $\left(\mathcal{P}_{w}^{N}\right)^{\sigma}=\mathcal{P}_{w^{\prime}}^{N}$, where $\left(\mathcal{P}_{w}^{N}\right)^{\sigma}$ is defined as in §2.7.

Proof. Since the points $t_{w}(a)$ are dense in $A_{w}$, we have $\left(A_{w}\right)^{\sigma}=A_{w^{\prime}}$. Clearly $\sigma$ sends the hyperplane sections of $A_{w}$ to those of $A_{w^{\prime}}$; also $t_{w}\left(u_{i}\right)^{\sigma}=t_{w^{\prime}}\left(u_{i}\right)$. Now, from (7.5) we obtain $\iota_{w}(c)^{\sigma} t_{w}(a)^{\sigma}=t_{w}(c a)^{\sigma}=t_{w^{\prime}}(c \lambda(a))=\iota_{w^{\prime}}(c) t_{w^{\prime}}(\lambda(a))=$ $\iota_{w^{\prime}}(c) t_{w}(a)^{\sigma}$ for every $c \in \mathfrak{r}$ and $a \in K_{2 n}^{1}$; thus $\iota_{w}(c)^{\sigma}=\iota_{w^{\prime}}(c)$. This completes the proof.
7.9. Since $\mathfrak{K} \subset \mathcal{A}_{0}=\mathbf{C} \mathfrak{F}$, we can find a countable subfield $k$ of $\mathbf{C}$ containing $k_{0}$ such that $\mathfrak{K} \subset k \mathfrak{F}$. Replacing $k$ by its algebraic closure in $\mathbf{C}$ if necessary, we assume that $k$ is algebraically closed. Now let $z_{0}$ be a generic point for the elements of $\mathfrak{F}$ over $k$. Take and fix a positive integer $N$. Since $\mathcal{P}_{z_{0}}^{N}$ is defined over $\mathfrak{K}\left[z_{0}\right]$, we can find elements $g_{1}, \ldots, g_{m}$ of $\mathfrak{K}$ such that $\mathcal{P}_{z_{0}}^{N}$ is defined over $\mathbf{Q}\left(g_{1}\left(z_{0}\right), \ldots, g_{m}\left(z_{0}\right)\right)$. Since the coordinates of $t_{z_{0}}(a)$ are algebraic over this field, $\mathfrak{K}\left[z_{0}\right]$ is algebraic over $\mathbf{Q}\left(g_{1}\left(z_{0}\right), \ldots, g_{m}\left(z_{0}\right)\right)$. Therefore $\mathfrak{K}$ is algebraic over $\mathbf{Q}\left(g_{1}, \ldots, g_{m}\right)$.

Let $V$ be the affine locus of the point $\left(g_{1}\left(z_{0}\right), \ldots, g_{m}\left(z_{0}\right)\right)$ over $k$. Take a multiple $M$ of $N$ so that the $g_{i}$ belong to $k \mathfrak{F}_{M}$. Since $z_{0}$ is generic for $\mathfrak{F}$ over $k, \varphi_{M}\left(z_{0}\right)$ is a generic point of $V_{M}$ over $k$, and we can define a $k$-rational map $q: V_{M} \rightarrow V$ by $q\left(\varphi_{M}\left(z_{0}\right)\right)=g\left(z_{0}\right)$, where we write $g=\left(g_{1},, \ldots, g_{m}\right)$.

Let us now prove that

$$
\begin{equation*}
k\left(\varphi_{N}\left(z_{0}\right)\right) \subset k\left(g\left(z_{0}\right)\right) \tag{7.9}
\end{equation*}
$$

Let $\sigma \in \operatorname{Aut}\left(\mathbf{C} / k\left(g\left(z_{0}\right)\right)\right)$. Since $V_{M}$ is defined over $k$, the point $\varphi_{M}\left(z_{0}\right)^{\sigma}$ belongs to $V_{M}$, and so $\varphi_{M}\left(z_{0}\right)^{\sigma}=\varphi_{M}\left(z_{1}\right)$ with $z_{1} \in \mathcal{H}$. This means that $z_{1}$ is generic for $\mathfrak{F}_{M}$ over $k$, and for $\mathfrak{F}$ over $k$ as well, since $\mathfrak{F}$ is algebraic over $\mathfrak{F}_{M}$. Thus $f\left(z_{0}\right) \mapsto f\left(z_{1}\right)$ for $f \in k \mathfrak{F}$ is an isomorphism; denote it by $\tau$. Then $g\left(z_{0}\right)^{\tau}=$ $q\left(\varphi_{M}\left(z_{0}\right)\right)^{\top}=q\left(\varphi_{M}\left(z_{1}\right)\right)=q\left(\varphi_{M}\left(z_{0}\right)\right)^{\sigma}=g\left(z_{0}\right)^{\sigma}=g\left(z_{0}\right)$. Since $\mathcal{P}_{z_{0}}^{N}$ is defined over $k\left(g\left(z_{0}\right)\right)$, we have $\left(\mathcal{P}_{z_{0}}^{N}\right)^{\tau}=\mathcal{P}_{z_{0}}^{N}$. On the other hand, $t_{z_{0}}(a)^{\tau}=t_{z_{1}}(a)$ for every $a \in K_{2 n}^{1}$, so that $\left(\mathcal{P}_{z_{0}}^{N}\right)^{\tau}=\mathcal{P}_{z_{1}}^{N_{0}}$ by Lemma 7.8. Thus $\mathcal{P}_{z_{0}}^{N}=\mathcal{P}_{z_{1}}^{N}$, and hence $\varphi_{N}\left(z_{0}\right)=\varphi_{N}\left(z_{1}\right)$. Since $\varphi_{N}\left(z_{1}\right)=p_{N}^{M}\left(\varphi_{M}\left(z_{1}\right)\right)=p_{N}^{M}\left(\varphi_{M}\left(z_{0}\right)\right)^{\sigma}=\varphi_{N}\left(z_{0}\right)^{\sigma}$, we obtain $\varphi_{N}\left(z_{0}\right)=\varphi_{N}\left(z_{0}\right)^{\sigma}$, which proves (7.9).
7.10. Theorem. (1) $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ is generated over $\mathbf{Q}$ by all the quotients of the form $\theta_{K}(0, r z ; \lambda) / \theta_{K}\left(0, r z ; \lambda^{\prime}\right)$ with $\mathbf{Q}_{\mathrm{ab}}$-valued $\lambda, \lambda^{\prime}$ in $\mathcal{S}\left(K_{\mathrm{h}}^{n}\right)$ and $0 \ll r \in F$.
(2) $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)=\mathfrak{K}$.
(3) $\mathcal{A}_{0}(\Phi)=\Phi \mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ for every subfield $\Phi$ of $\mathbf{C}$ containing $\mathbf{Q}_{\mathrm{ab}}$; in particular, $\mathcal{A}_{0}=\mathbf{C} \mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$.
(4) $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ is stable under the map $f \mapsto f \circ \alpha$ for every $\alpha \in \widetilde{G}_{+}$.

Proof. From (7.9) we obtain $\mathfrak{F}^{N} \subset k \mathfrak{K}$, and consequently $\mathcal{A}_{0}=\mathbf{C} \mathfrak{F}=\mathbf{C} \mathfrak{K}$. Since $\mathbf{Q}_{\mathrm{ab}} \subset \mathfrak{K} \subset \mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, we obtain $\mathcal{A}_{0}=\mathbf{C} \mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$. Now, given a subfield $\Phi$ of $\mathbf{C}$ containing $\mathbf{Q}_{\mathrm{ab}}$, we have $\Phi \mathfrak{K} \subset \Phi \mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right) \subset \mathcal{A}_{0}(\Phi)$ and $\mathbf{C} \mathcal{A}_{0}(\Phi) \subset \mathcal{A}_{0}=\mathbf{C} \mathfrak{K}=$
$\mathbf{C} \Phi \mathfrak{K}$, so that $\mathbf{C} \mathcal{A}_{0}(\Phi)=\mathbf{C} \Phi \mathfrak{K}$. By Lemma $7.4(1), \mathcal{A}_{0}(\Phi)$ and $\mathbf{C}$ are linearly disjoint over $\Phi$. Therefore, by Lemma 7.2, $\Phi \mathfrak{K}=\mathcal{A}_{0}(\Phi)$. Taking $\Phi=\mathbf{Q}_{\mathrm{ab}}$, we obtain (2), and hence (3) as well. As for (1), Proposition 6.14 shows that any quotient of the form (7.8) is of the form described in (1). Now let $f$ and $g$ denote the numerator and the denominator of the quotient of (1). In Case SP, by Propositions 6.14 and 6.15 , both $f g$ and $g^{2}$ belong to $\mathcal{M}_{1}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, and hence $f / g \in \mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$. Clearly the same conclusion is true in Case UT. This proves (1). As for (4), by Lemma 7.5 it is sufficient to prove the cases $\alpha=\eta_{n}$ and $\alpha \in P \cap \widetilde{G}_{+}$. Since $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)=\mathfrak{K}$, our task is to show that $\alpha$ sends a quotient of the form (7.8) into $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$. If $\alpha=\eta_{n}$, we have $r \cdot \eta(z)=\eta\left(r^{-1} z\right)$. By Theorem 6.9 (2) or Theorem 6.11 (1) we have

$$
\begin{aligned}
\theta_{K}\left(0, r \cdot \eta(z) ; \lambda_{1}\right) / \theta_{K}\left(0, r \cdot \eta(z) ; l_{2}\right) & =\theta_{K}\left(0, \eta\left(r^{-1} z\right) ; \lambda_{1}\right) / \theta_{K}\left(0, \eta\left(r^{-1} z\right) ; \lambda_{2}\right) \\
& =\theta_{K}\left(0, r^{-1} z ; \lambda_{1}^{\prime}\right) / \theta_{K}\left(0, r^{-1} z ; \lambda_{2}^{\prime}\right),
\end{aligned}
$$

where $\lambda_{i}^{\prime}$ can be obtained from $\lambda_{i}$ by (6.33) or Theorem 6.11 (4). We easily see that $\lambda_{i}^{\prime}$ is $\mathbf{Q}_{\mathrm{ab}}$-valued if $\lambda_{i}$ is so. Thus $\eta$ sends $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ into itself. The action of $P \cap \widetilde{G}_{+}$will be discussed in the proof of the following theorem.
7.11. Theorem. Let $\kappa \in \mathbf{Z}$ in Case UT and $\kappa \in 2^{-1} \mathbf{Z}$ in Case SP. Given $\alpha \in \widetilde{G}_{+}$, let $q(z)=\left(j_{\alpha}(z)^{\mathbf{a}}\right)^{\kappa}$, where we take any branch of the square root of $j_{\alpha}(z)^{2 \kappa \mathbf{a}}$ to be $q(z)$ if $\kappa \notin \mathbf{Z}$. Let $\Phi$ be a subfield of $\mathbf{C}$ containing $\mathbf{Q}_{\mathrm{ab}}$. If $\kappa \neq 0$ in Case UT, then assume either $\operatorname{det}(\alpha)=\nu(\alpha)^{n}$ or $\Phi$ contains the reflex field $K^{\prime}$ of $\S 1.12$ defined for $(K, \tau)$ of $\S 3.5$. Then $q(z)^{-1} f(\alpha z) \in \mathcal{M}_{\kappa \mathbf{a}}(\Phi)$ for every $f \in \mathcal{M}_{\kappa \mathbf{a}}(\Phi)$. Moreover, in Case UT, $f \|_{\kappa \mathbf{b}} \alpha \in \mathcal{M}_{\kappa \mathbf{b}}(\Phi)$ for every $f \in \mathcal{M}_{\kappa \mathbf{b}}(\Phi)$, every $\alpha \in \widetilde{G}_{+}$, and every subfield $\Phi$ of $\mathbf{C}$ containing $\mathbf{Q}_{\mathrm{ab}}$.

Proof. We first assume $K^{\prime} \subset \Phi$ in Case UT. By Lemma 7.5 it is sufficient to prove the cases $\alpha=\eta_{n}$ and $\alpha \in P \cap \widetilde{G}_{+}$. (Notice that $\eta_{n}^{-1}=-\eta_{n} \in \eta_{n}\left(P \cap \widetilde{G}_{+}\right)$.) Let $\alpha=\left[\begin{array}{ll}a & b \\ 0 & d\end{array}\right] \in P \cap \widetilde{G}_{+}, s=\nu(\alpha)$, and $f(z)=\sum_{h} c(h) \mathbf{e}_{\mathbf{a}}^{n}(h z) \in \mathcal{M}_{\kappa \mathbf{a}}(\Phi)$. Then $a^{*} d=s 1_{n}, q(z)=\prod_{v \in \mathbf{a}} \operatorname{det}(d)_{v}^{\kappa} \in K^{\prime}$, and

$$
q(z)^{-1} f(\alpha z)=q(z)^{-1} \sum_{h} c(h) \mathbf{e}_{\mathbf{a}}^{n}\left(s^{-1} h b a^{*}\right) \mathbf{e}_{\mathbf{a}}^{n}\left(s^{-1} a^{*} h a z\right) .
$$

Since $b a^{*} \in S$ by (1.13), we see that $\operatorname{tr}\left(s^{-1} h b a^{*}\right) \in F$, and hence $\mathbf{e}_{\mathbf{a}}^{n}\left(s^{-1} h b a^{*}\right)$ is a root of unity; thus $q(z)^{-1} f(\alpha z) \in \mathcal{M}_{\kappa \mathbf{a}}(\Phi)$. In Case SP the same reasoning is valid for any $\Phi$ containing $\mathbf{Q}_{\mathrm{ab}}$.

Returning to the proof of Theorem 7.10 (4), we apply the same technique to $f=\theta_{K}(0, r z ; \lambda)$ with a $\mathbf{Q}_{\mathrm{ab}}$-valued $\lambda$. Then for $\alpha \in P \cap \widetilde{G}_{+}$the above reasoning shows that $f \circ \alpha$ belongs to $\mathcal{M}_{\mathbf{a} / 2}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ or $\mathcal{M}_{\mathrm{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, and so $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ is stable under $P \cap \widetilde{G}_{+}$. This completes the proof of Theorem 7.10.

Next we consider the action of $\eta$ on $\mathcal{M}_{\kappa \mathbf{a}}(\Phi)$. Given $f \in \mathcal{M}_{\kappa \mathbf{a}}(\Phi)$, put $m=$ $2 \kappa /[K: F]$ and $t=f g^{-m}$ with a nonzero function $g$ of the form $g(z)=\theta_{K}(0, z ; \lambda)$ with a $\mathbf{Q}_{\mathrm{ab}}$-valued $\lambda$. Then $g^{m} \in \mathcal{M}_{\kappa \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ and $t \in \mathcal{A}_{0}(\Phi)$. By Theorem 7.10 (3), (4), $\mathcal{A}_{0}(\Phi)$ is stable under $\widetilde{G}_{+}$. In Case SP we have $q(z)=\zeta \cdot h_{\eta}(z)^{m}$ with a root of unity $\zeta$, and $q(z)^{-1} f(\eta z)=\zeta^{-1}(g \| \eta)^{m}(t \circ \eta)$, which belongs to $\mathcal{A}_{\kappa \mathbf{a}}(\Phi)$, since $t \circ \eta \in \mathcal{A}_{0}(\Phi)$ and $(g \| \eta)^{m} \in \mathcal{M}_{\kappa \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$. The last fact follows from the behavior of $\theta_{K}$ under $\eta$ as discussed at the end of the proof of Theorem 7.10 (cf. Proposition 6.15). By (5.30), $q(z)^{-1} f(\eta z) \in \mathcal{M}_{\kappa \mathbf{a}}(\Phi)$. Next, take $\alpha \in \widetilde{G}_{+}$such that $\operatorname{det}(\alpha)=\nu(\alpha)^{n}$ in Case UT. Given such an $\alpha$, put $\tau=\operatorname{diag}\left[\nu(\alpha) 1_{n}, 1_{n}\right]$
and $\beta=\tau^{-1} \alpha$. Then $\beta \in G_{1}$. By Lemma 7.5, $G_{1}$ is generated by $\eta$ and $G_{1} \cap P$. If $\gamma \in G_{1} \cap P$ or $\gamma=\tau$, then $\operatorname{det}\left(d_{\gamma}\right) \in F$, so that $j_{\gamma}(z)^{\mathbf{a}} \in \mathbf{Q}$. Therefore our reasoning is valid for the group of elements $\alpha$ satisfying $\operatorname{det}(\alpha)=\nu(\alpha)^{n}$ without the condition $K^{\prime} \subset \Phi$. Finally take $\mathcal{M}_{\kappa \mathbf{b}}$ in Case UT. Since $j_{\gamma}^{\mathbf{b}} \in \mathbf{Q}$ for $\gamma \in P \cap \widetilde{G}_{+}$, we do not have to assume that $K^{\prime} \subset \Phi$. Thus we obtain the last assertion.
7.12. Lemma. If $\left(\mathcal{P}_{w}^{N}\right)^{\sigma}$ is isomorphic to $\mathcal{P}_{z}^{N}$ for $\sigma \in \operatorname{Aut}(\mathbf{C})$ and some $z, w \in$ $\mathcal{H}$, then $\mathbf{e}(1 / N)^{\sigma}=\mathbf{e}(1 / N)$.

Proof. Let $s$ be the positive integer such that $\left\{\operatorname{Tr}_{K / \mathbf{Q}}\left(x \eta_{n} y^{*}\right) \mid x, y \in L\right\}=$ $s \mathbf{Z}$. Then there exists a divisor $X_{z}$ on $A_{z}$ that determines the Riemann form $E_{z}$ of (4.25) with $\mathcal{T}=s^{-1} \eta_{n}$. Then $X_{z}$ is a basic polar divisor of $\mathcal{P}_{z}$, and by (2.10) we have

$$
\begin{equation*}
\zeta_{X_{z}}\left(t_{z}(a), t_{z}(b)\right)=\mathbf{e}\left(m \cdot E_{z}\left(p_{z}(a), p_{z}(b)\right)\right)=\mathbf{e}\left((m / s) \operatorname{Tr}_{K / \mathbf{Q}}\left(a \eta_{n} b^{*}\right)\right) \tag{7.10}
\end{equation*}
$$

for every $a, b \in m^{-1} L, 0<m \in \mathbf{Z}$. Let $\varepsilon$ be an isomorphism of $\left(\mathcal{P}_{w}^{N}\right)^{\sigma}$ onto $\mathcal{P}_{z}^{N}$. Then $\varepsilon t_{w}(a)^{\sigma}=t_{z}(a)$ for every $a \in N^{-1} L$. Let $Y=\varepsilon\left(X_{w}^{\sigma}\right)$. Since $X_{w}$ is a basic polar divisor of $\mathcal{P}_{w},\left(X_{w}\right)^{\sigma}$ is a basic polar divisor of $\left(\mathcal{P}_{w}\right)^{\sigma}$. (This follows from the characterization of a basic polar divisor in terms of algebraic equivalence mentioned in §2.3.) Therefore $Y$ is a basic polar divisor of $\mathcal{P}_{z}$, so that $Y$ determines $E_{z}$. Therefore for $a, b \in N^{-1} L$ we have

$$
\begin{gathered}
\mathbf{e}\left((N / s) \operatorname{Tr}_{K / \mathbf{Q}}\left(a \eta_{n} b^{*}\right)\right)=\zeta_{Y}\left(t_{z}(a), t_{z}(b)\right)=\zeta_{X_{w}^{\sigma}}\left(t_{w}(a)^{\sigma}, t_{w}(b)^{\sigma}\right) \\
=\zeta_{X_{w}}\left(t_{w}(a), t_{w}(b)\right)^{\sigma}=\mathbf{e}\left((N / s) \operatorname{Tr}_{K / \mathbf{Q}}\left(a \eta_{n} b^{*}\right)\right)^{\sigma}
\end{gathered}
$$

by (2.11). We can take $a$ and $b$ in $N^{-1} L$ so that $\operatorname{Tr}_{K / Q}\left(a \eta_{n} b^{*}\right)=s N^{-2}$. Then we obtain $\mathbf{e}(1 / N)^{\sigma}=\mathbf{e}(1 / N)$ as expected.

## 8. Action of certain elements of $\widetilde{G}_{\text {A }}$ on $\mathfrak{K}$

8.1. It is well-known that $\mathbf{Q}^{\times} \mathbf{Q}_{\mathbf{a}+}^{\times}$is closed in $\mathbf{Q}_{\mathbf{A}}^{\times}$, and there is a canonical homomorphism of $\mathbf{Q}_{\mathbf{A}}^{\times}$onto $\operatorname{Gal}\left(\mathbf{Q}_{\mathrm{ab}} / \mathbf{Q}\right)$ with kernel $\mathbf{Q}^{\times} \mathbf{Q}_{\mathbf{a}+}^{\times}$. For $t \in \mathbf{Q}_{\mathbf{A}}^{\times}$we denote by $[t, \mathbf{Q}]$ the element of $\operatorname{Gal}\left(\mathbf{Q}_{\mathbf{a b}} / \mathbf{Q}\right)$ that is the image of $t$ under that homomorphism. To simplify our notation, let us now put

$$
\begin{equation*}
\mathbf{Z}_{\mathbf{h}}^{\times}=\prod_{p} \mathbf{Z}_{p}^{\times} \tag{8.1}
\end{equation*}
$$

where the product is taken over all rational primes $p$. Observing that $\mathbf{Q}_{\mathbf{A}}^{\times} / \mathbf{Q}^{\times} \mathbf{Q}_{\mathbf{a}+}^{\times}$ is isomorphic to $\mathbf{Z}_{\mathbf{h}}^{\times}$, we see that the map $t \rightarrow[t, \mathbf{Q}]$ gives an isomorphism of $\mathbf{Z}_{\mathbf{h}}^{\times}$ onto $\operatorname{Gal}\left(\mathbf{Q}_{\mathrm{ab}} / \mathbf{Q}\right)$. Now $\mathbf{Q} / \mathbf{Z}$ is canonically isomorphic to $\mathbf{Q}_{\mathbf{A}} /\left(\mathbf{Q}_{\mathbf{a}} \prod_{p} \mathbf{Z}_{p}\right)$, and we can let $\mathbf{Z}_{\mathbf{h}}^{\times}$act on the last group by multiplication. For $t \in \mathbf{Z}_{\mathbf{h}}^{\times}$and $x \in \mathbf{Q} / \mathbf{Z}$ denote by $t x$ the image of $x$ under $t$. Clearly $x \rightarrow \mathbf{e}_{\mathbf{h}}(x)=\prod_{p} \mathbf{e}_{p}(x)$ for $x \in \mathbf{Q}$ gives an isomorphism of $\mathbf{Q} / \mathbf{Z}$ onto the group of all roots of unity. (See $\S 1.6$ for the definition of $\mathbf{e}_{p}$. Notice also that $\mathbf{e}_{\mathbf{h}}(x)=\mathbf{e}(-x)$ for every $x \in \mathbf{Q}$.) Then we can easily show that

$$
\begin{equation*}
\mathbf{e}_{\mathbf{h}}(x)^{[t, \mathbf{Q}]}=\mathbf{e}_{\mathbf{h}}\left(t^{-1} x\right) \quad\left(t \in \mathbf{Z}_{\mathbf{h}}^{\times}, x \in \mathbf{Q} / \mathbf{Z}\right) \tag{8.2}
\end{equation*}
$$

In particular, given $t \in \mathbf{Z}_{\mathbf{h}}^{\times}$and a positive integer $N$, take a positive integer $r$ so that $r t_{p}-1 \in N \mathbf{Z}_{p}$ for every prime $p$. Then (8.2) means $\mathbf{e}(1 / N)^{[t, \mathbf{Q}]}=\mathbf{e}(r / N)$, which is the classical reciprocity-law in cyclotomic fields.
8.2. Our setting is the same as in Sections 6 and 7 ; see also $\S 5.6$. We retain the convention that $K=F$ and $\mathfrak{r}=\mathfrak{g}$ in Case SP. We use the symbols $\widetilde{G}, G$, and $G_{0}$ of (3.26), (3.27), and (3.29) in both cases. Put $G_{1}=G \cap S L_{2 n}(K)$ as in Lemma 7.5. Then $G_{0} \cap G=G_{1} ; G=G_{1}=S p(n, F)$ in Case SP. Define $\iota: \mathbf{Q}_{\mathbf{A}}^{\times} \rightarrow\left(G_{0}\right)_{\mathbf{A}}$ by

$$
\begin{equation*}
\iota(s)=\operatorname{diag}\left[1_{n}, s^{-1} 1_{n}\right] \quad\left(s \in \mathbf{Q}_{\mathbf{A}}^{\times}\right) \tag{8.3}
\end{equation*}
$$

We now take the $\mathfrak{r}$-lattice $L$ of (7.4b) to be $\mathfrak{r}_{2 n}^{1}$ (Actually our treatment is applicable to a more general type of $L, L=\mathfrak{r}_{2 n}^{1} \xi$, for example, with any diagonal $\xi$ in $G_{\mathbf{h}}$ such that $\operatorname{Tr}_{K / \mathbf{Q}}\left(x \eta_{n} y^{*}\right) \in \mathbf{Z}$ for every $x, y \in L$.) We then put

$$
\begin{align*}
& U^{1}=\left\{x \in\left(G_{0}\right)_{\mathbf{A}+} \mid L x=L\right\}  \tag{8.4}\\
& U^{N}=\left\{x \in U^{1} \mid L_{v}\left(x_{v}-1\right) \subset N L_{v} \text { for every } v \in \mathbf{h}\right\} \quad(0<N \in \mathbf{Z})  \tag{8.5}\\
& T^{N}=\iota\left(\mathbf{Z}_{\mathbf{h}}^{\times}\right) U^{N} \tag{8.6}
\end{align*}
$$

Clearly $U^{1}$ is a subgroup of $\left(G_{0}\right)_{\mathbf{A}+}$ and $U^{N}$ is a normal aubgroup of $U^{1}$; since $\iota\left(\mathbf{Z}_{\mathbf{h}}^{\times}\right) \subset U^{1}, T^{N}$ is a subgroup of $U^{1}$ and $T^{N}=U^{N} \iota\left(\mathbf{Z}_{\mathbf{h}}^{\times}\right)$. We also employ $\Gamma^{N}$ of (7.6).
8.3. Lemma. (1) $\left(G_{0}\right)_{\mathbf{A}+}=G_{0+} T^{N}=T^{N} G_{0+}$ for every $N$.
(2) $\widetilde{G} \cap T^{N} \widetilde{G}_{\mathbf{a}+}=\Gamma^{N} \cap T^{N}=U^{N} \cap G_{1}$ for every $N$.
(3) $\Gamma^{N} \subset U^{N} \cap G_{1}$ for $N>2$ in Case UT and for every $N$ in Case $S P$.
(4) Given $x \in T^{N} \Gamma^{N}$ and a multiple $M$ of $N$, there exists an element $\gamma$ of $\Gamma^{N}$ and an element $y$ of $U^{M}$ such that $x=\iota(r) y \gamma$, where $r=\nu(x)^{-1}$.

Proof. Let $x \in\left(G_{0}\right)_{\mathbf{A}+\text {. Then }} 0 \ll \nu(x) \in \mathbf{Q}_{\mathbf{A}}^{\times}$, and hence $\nu(x)=a b c$ with $0<a \in \mathbf{Q}^{\times}, b \in \mathbf{Z}_{\mathbf{h}}^{\times}$, and $0<c \in \mathbf{Q}_{\mathbf{a}}^{\times}$. Put $y=\iota(a) x \iota(b c)$. Then $\nu(y)=1$ and $y \in\left(G_{0}\right)_{\mathbf{A}}$, so that $y \in\left(G_{1}\right)_{\mathbf{A}}$. By strong approximation in $G_{1}$ we have $\left(G_{1}\right)_{\mathbf{A}} \subset$ $G_{1} U^{N}$ for every $N$, and hence $x=\iota(a)^{-1} y \iota(b c)^{-1} \in G_{0+} G_{1} U^{N} \iota(b)^{-1} \subset G_{0+} T^{N}$, from which we can easily derive (1). To prove (2), let $\gamma \in \widetilde{G} \cap T^{N} \widetilde{G}_{\mathbf{a}_{+}}$. Then $\gamma \in \iota(s) x \widetilde{G}_{\mathbf{a}+}$ with $s \in \mathbf{Z}_{\mathbf{h}}^{\times}$and $x \in U^{N}$. We see that $\nu(\gamma)$ is a totally positive unit belonging to $\mathbf{Q}$, and so $\nu(\gamma)=1$. Similarly $\operatorname{det}(\gamma)=1$. Thus $s=\nu(x)_{\mathbf{h}}$. Since $x \in U^{N}$, we have $s-1 \prec N \mathbf{Z}$, so that $\gamma-1 \prec N$ r. From these facts we easily obtain (2). By (4.34), $\Gamma^{N} \subset G_{1}$ for $N>2$ in Case UT and for every $N$. in Case SP, which implies (3). Clearly it is sufficient to prove (4) for $x \in T^{N}$. Given $x \in T^{N}$, put $r=\nu(x)^{-1}$ and $z=\iota(r)^{-1} x$. Then $\nu(z)=1$; also $z \in T^{M} G_{0}$ by (1), and hence $z=\iota(s) w \gamma$ with $s \in \mathbf{Z}_{\mathbf{h}}^{\times}, w \in U^{M}$, and $\gamma \in G_{0}$. Then $\gamma \in G_{0} \cap T^{N} \subset \Gamma^{N}$ by (2), and hence $\nu(\gamma)=1$. Thus $s=\nu(w)$, and $s-1 \prec M \mathbf{Z}$ since $w \in U^{M}$. Put $y=\iota(s) w$; then $y \in U^{M}$ and $x=\iota(r) y \gamma$. This proves (4).
8.4. We are going to define the action of a certain subgroup of $\widetilde{G}_{\mathbf{A}}$ on $\mathfrak{K}$. First, in view of Theorem 7.10 (4), the map $f \rightarrow f \circ \alpha$ is clearly an automorphism of $\mathfrak{K}$ for every $\alpha \in \widetilde{G}_{+}$.

Observe that $K_{2 n}^{1} / L$ is canonically isomorphic to $\left(K_{2 n}^{1}\right)_{\mathbf{A}} /\left[\left(K_{2 n}^{1}\right)_{\mathbf{a}} \prod_{v \in \mathbf{h}} L_{v}\right]$. Therefore we can define $t(a, z)$ for $a \in\left(K_{2 n}^{1}\right)_{\mathbf{A}}$ by putting $t(a, z)=t(b, z)$ with $b \in K_{2 n}^{1}$ such that $a_{v}-b \in L_{v}$ for every $v \in \mathbf{h}$. In particular, if $a \in K_{2 n}^{1}$ and $x \in T^{1}$, then $a x(\bmod L)$ and $t(a x, w)$ are meaningful, and they depend only on $a$ modulo $L$. Notice that $t(a x, z)=t(a, z)$ for every $a \in K_{2 n}^{1}$ if and only if $x \in \widetilde{G}_{\mathbf{a}+}$.

We now consider the action of $\operatorname{Aut}(\mathbf{C})$ on the field of quotients of the formal series defined in §5.9. In particular we let $\operatorname{Gal}\left(\mathbf{Q}_{\mathrm{ab}} / \mathbf{Q}\right)$ act on $\mathfrak{K}=\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$. For the moment $f^{\sigma}$ for $f \in \mathfrak{K}$ and $\sigma \in \operatorname{Gal}\left(\mathbf{Q}_{\mathrm{ab}} / \mathbf{Q}\right)$ is just the quotient of two formal
series, but the following lemma will show that $f^{\sigma} \in \mathfrak{K}$. For a fixed $a \in K_{2 n}^{1}$ the quotients of the projective coordinates of $t(a, z)$ as functions of $z$ belong to $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, and so we can let $\operatorname{Gal}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ act on $t(a, z)$.
8.5. Lemma. (1) We have

$$
\begin{equation*}
t(a, z)^{[s, \mathbf{Q}]}=t(a \iota(s), z) \quad \text { for every } \quad s \in \mathbf{Z}_{\mathbf{h}}^{\times} \quad \text { and } \quad \text { every } \quad a \in K_{2 n}^{1} . \tag{8.7}
\end{equation*}
$$

Consequently $\mathfrak{K}$ is stable under $\operatorname{Gal}\left(\mathbf{Q}_{\mathrm{ab}} / \mathbf{Q}\right)$.
(2) Every element of $\mathfrak{K}$ is of the form $q / r$ with $q \in \mathcal{M}_{\kappa \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ and $0 \neq r \in$ $\mathcal{M}_{\kappa \mathbf{a}}(\mathbf{Q})$, for some positive integer $\kappa$, such that $q^{\sigma} \in \mathcal{M}_{\kappa \mathbf{a}}\left(\mathbf{Q}_{\mathbf{a b}}\right)$ for every $\sigma \in$ $\operatorname{Gal}\left(\mathbf{Q}_{\mathrm{ab}} / \mathbf{Q}\right)$.
(3) For every subfield $D$ of $\mathrm{Q}_{\mathrm{ab}}$ we have

$$
\mathcal{A}_{0}(D)=\left\{f \in \mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right) \mid f^{\sigma}=f \quad \text { for every } \sigma \in \operatorname{Gal}\left(\mathbf{Q}_{\mathrm{ab}} / D\right)\right\} .
$$

Proof. We prove (1) only in Case UT. Case SP is similar and simpler. Take $\Lambda$ as in Theorem 6.12 and $\S 7.6$. Let $a=\left({ }^{t} b,{ }^{t} c\right) \in K_{2 n}^{1}$ with $b, c \in N^{-1} \mathfrak{r}_{1}^{n}, 0<N \in \mathbf{Z}$. By (4.24), $p_{z}(a)=\left(z \bar{b}+\bar{c},{ }^{t} z b+c\right)$, and hence, by Proposition 6.14, for $\ell_{1}, \ell_{2} \in \Lambda$ we have

$$
\theta_{K}\left(p_{z}(a), p^{-1} z ; \ell_{1}\right) / \theta_{K}\left(p_{z}(a), p^{-1} z ; \ell_{2}\right)=\theta_{K}\left(0, p^{-1} z ; \ell_{1}^{\prime}\right) / \theta_{K}\left(0, p^{-1} z ; \ell_{2}^{\prime}\right)
$$

where $\ell_{i}^{\prime}(h)=\ell_{i}(h-p b) \mathbf{e}\left(\operatorname{Tr}_{K / \mathbf{Q}}\left(c^{*}(h-p b)\right)\right)$ for $h \in K^{n}$. Take a multiple $M^{-}$of $N$ so that $\ell(h)=0$ for every $\ell \in \Lambda$ if $M h \notin \mathfrak{r}^{n}$. Given $s \in \mathbf{Z}_{\mathbf{h}}^{\times}$, take $c_{1} \in K^{n}$ so that $\left(c_{1}-s^{-1} c\right)_{v} \in M \mathfrak{r}_{v}^{n}$ for every $v \in \mathbf{h}$. Since $\ell_{i}$ is $\mathbf{Z}$-valued, by (8.2), $\sigma=[s, \mathbf{Q}]$ sends the last quotient to $\theta_{K}\left(0 ; p^{-1} z ; \lambda_{1}\right) / \theta_{K}\left(0 ; p^{-1} z ; \lambda_{2}\right)$ with $\lambda_{i}(h)=$ $\ell_{i}(h-p b) \mathbf{e}\left(\operatorname{Tr}_{K / \mathbf{Q}}\left(c_{1}^{*}(h-p b)\right)\right)$. This means that $t_{z}(a)^{\sigma}=t_{z}\left(\left({ }^{t} b,{ }^{t} c_{1}\right)\right)$, which proves (1).

Next, from this argument we can easily derive that every element of $\mathfrak{K}$ is of the form $g_{1} / g$ with $g, g_{1} \in \mathcal{M}_{\kappa \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ for some $\kappa$ such that $g_{1}^{\sigma}, g^{\sigma} \in \mathcal{M}_{\kappa \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ for every $\sigma \in \operatorname{Gal}\left(\mathbf{Q}_{\mathrm{ab}} / \mathbf{Q}\right)$. (In Case SP, multiplying both denominator and numerator of (7.8) by the denominator, and employing Proposition 6.15, we can avoid halfintegral weights.) Also $g \in \mathcal{M}_{\kappa \mathbf{a}}(\Phi)$ with a finite extension $\Phi$ of $\mathbf{Q}$ contained in $\mathbf{Q}_{\mathrm{ab}}$. Then taking $r$ to be the product of $g^{\sigma}$ for all $\sigma \in \operatorname{Gal}(\Phi / \mathbf{Q})$ and putting $q=g_{1} r / g$, we obtain (2). Assertion (3) follows immediately from (2).

We now consider the following condition on $\Gamma^{N}$ :

$$
\begin{equation*}
t(a \gamma, z)=t(a, \gamma z) \quad \text { for every } \quad a \in K_{2 n}^{1} \text { and } \gamma \in \Gamma^{N} \tag{8.8}
\end{equation*}
$$

Let $\Gamma_{0}$ be the group as in Theorem 6.12 which we fixed in $\S 7.6$. Then, by (7.3), condition (8.8) is satified if $\Gamma^{N} \subset \Gamma_{0}$, which is so for any multiple $N$ of a suitably chosen positive integer.
8.6. Lemma. Under condition (8.8), given $x \in T^{N}$, there is a unique automorphism $\tau(x)$ of $\mathfrak{K}$ such that $t(a, z)^{\tau(x)}=t(a x, z)$ for every $a \in K_{2 n}^{1}$. Moreover, $\tau(x)=\left[\nu(x)^{-1}, \mathbf{Q}\right]$ on $\mathbf{Q}_{\mathrm{ab}}, \tau(x y)=\tau(x) \tau(y)$, and $f^{\tau(x)}=f \circ x$ for $f \in \mathfrak{K}$ if $x \in \Gamma^{N}$.

Proof. The uniqueness is obvious, since $\mathfrak{K}$ is generated by the affine coordinates of $t(a, z)$. Also the statement is true if $x \in \Gamma^{N}$, since (8.8) shows that the desired automorhism can be given by $f \rightarrow f \circ x$, which is the last assertion. Take a multiple $M$ of $N$. Let $\Phi_{M}$ denote the field generated over $\mathbf{Q}$ by the affine coordinates of $t(a, w)$ for all $a \in M^{-1} L$. Now, given $x \in T^{N}$, by Lemma 8.3 (4) we can put $x=\iota(r) y \gamma$ with $r=\nu(x)^{-1}, y \in U^{M}$, and $\gamma \in \Gamma^{N}$. Then for $a \in M^{-1} L$ we have
$t(a x, z)=t(a \iota(r) y \gamma, z)=t(a \iota(r) y, \gamma z)=t(a \iota(r), \gamma z)=\left(t(a, z)^{[r . \mathbf{Q}]}\right) \circ \gamma$
by (8.7) and (8.8). This means that there is an automorphism $\sigma$ of $\Phi_{M}$ that sends $t(a, z)$ to $t(a x, z)$ for every $a \in M^{-1} L$; moreover, $f^{\sigma}=\left(f^{[r, \mathbf{Q}]}\right) \circ \gamma$. Thus $\sigma=[r, \mathbf{Q}]$ on $\mathbf{Q}_{\mathrm{ab}} \cap \Phi_{M}$. Taking all multiples $M$ of $N$, we obtain the desired automorphism $\tau(x)$. The equality $\tau(x y)=\tau(x) \tau(y)$ is obvious from the definition.
8.7. Lemma. Under (8.8) let $\alpha, \beta \in \widetilde{G}_{+}$and $x, y \in T^{N}$; suppose $y \beta \in \alpha x \widetilde{G}_{\mathbf{a}+}$. Then $(f \circ \alpha)^{\tau(x)}=f^{\tau(y)} \circ \beta$ for every $f \in \mathfrak{K}$.

Proof. Put $w=x^{-1} \alpha^{-1} y \beta$. Then $w \in \widetilde{G}_{\mathbf{a}+} \cap\left(G_{0}\right)_{\mathbf{A}} \subset T^{N}$ and $\tau(x)=\tau(x w)$. Therefore, changing $x$ for $x w$, we may assume that $\alpha x=y \beta$. Changing also $\alpha$ and $\beta$ for $c \alpha$ and $c \beta$ with a suitable $c \in \mathbf{Q}$, we may assume that $\alpha^{-1}, \beta^{-1} \prec \mathfrak{r}$. Take a generic point $z_{0}$ for $\mathfrak{F}$ over $k$ as in $\S 7.9$. Since the $\operatorname{map} f \mapsto f\left(z_{0}\right)$ is an isomorphism, we can define an automorphism $\sigma$ of $\mathfrak{K}\left[z_{0}\right]$ by $f\left(z_{0}\right)^{\sigma}=f^{\tau(x)}\left(z_{0}\right)$ for $f \in \mathfrak{K}$. Then our assertion is equivalent to the equality $t\left(a, \alpha z_{0}\right)^{\sigma}=t\left(a y, \beta z_{0}\right)$. To simplify our notation, put $\alpha=\alpha_{1}, \beta=\alpha_{2}$, and $z_{\nu}=\alpha_{\nu} z_{0}$. Observe that $\mathfrak{K}\left[z_{\nu}\right]=\mathfrak{K}\left[z_{0}\right]$ and $A\left(z_{0}\right)^{\sigma}=A\left(z_{0}\right)$ since $t\left(a, z_{0}\right)^{\sigma}=t\left(a x, z_{0}\right)$ by our definition. Now we consider the following commutative diagram:


Here $\Lambda_{\nu}={ }^{t} M\left(\alpha_{\nu}, z_{0}\right)^{-1}$ with the symbol $M$ of (4.29). From (4.31) we obtain $\Lambda_{\nu} p\left(a, z_{0}\right)=p\left(a \alpha_{\nu}^{-1}, z_{\nu}\right)$, which gives the leftmost part of the diagram; $\Theta_{z}(u)=$ $\Theta_{K}(u, z)$ with $\Theta_{K}$ of Theorem $6.12 ; \lambda_{\nu}$ is an isogeny determined by $\Lambda_{\nu}$ with the property that $\lambda_{\nu} t\left(a, z_{0}\right)=t\left(a \alpha_{\nu}^{-1}, z_{\nu}\right)$ and $\operatorname{Ker}\left(\lambda_{\nu}\right)=t\left(L \alpha_{\nu}, z_{0}\right)$. Notice that $\lambda_{\nu} t\left(a, z_{0}\right)=t\left(a \alpha_{\nu}^{-1}, z_{\nu}\right)$ holds even for $a \in\left(K_{2 n}^{1}\right)_{\mathbf{A}}$. Since $t\left(a, z_{0}\right)^{\sigma}=t\left(a x, z_{0}\right)$, we have $\operatorname{Ker}\left(\lambda_{1}\right)^{\sigma}=t\left(L \alpha_{1} x, z_{0}\right)=t\left(L y \alpha_{2}, z_{0}\right)=t\left(L \alpha_{2}, z_{0}\right)=\operatorname{Ker}\left(\lambda_{2}\right)$, and hence there exists an isomorphism $\varepsilon$ of $A\left(z_{2}\right)$ onto $A\left(z_{1}\right)^{\sigma}$ such that $\lambda_{1}^{\sigma}=\varepsilon \lambda_{2}$. Observe that $A\left(z_{\nu}\right)$ and $\lambda_{\nu}$ are rational over $\mathfrak{K}\left[z_{0}\right]$. Now for every $a \in K_{2 n}^{1}$, we have $t\left(a \alpha_{1}^{-1}, z_{1}\right)^{\sigma}=\lambda_{1}^{\sigma} t\left(a, z_{0}\right)^{\sigma}=\lambda_{1}^{\sigma} t\left(a x, z_{0}\right)=\varepsilon \lambda_{2} t\left(a x, z_{0}\right)=\varepsilon t\left(a x \alpha_{2}^{-1}, z_{2}\right)=$ $\varepsilon t\left(a \alpha_{1}^{-1} y, z_{2}\right)$, which shows that

$$
\begin{equation*}
t\left(a, z_{1}\right)^{\sigma}=\varepsilon t\left(a y, z_{2}\right) \tag{*}
\end{equation*}
$$

Put $\mathcal{P}(z)=\left(A_{z}, \mathcal{C}_{z}, \iota_{z}\right)$. Then $\varepsilon$ gives an isomorphism of $\mathcal{P}\left(z_{2}\right)$ onto $\mathcal{P}\left(z_{1}\right)^{\sigma}$, which will be proven at the end of the proof. Taking $\left(y^{-1}, z_{2}\right)$ in place of $\left(x, z_{0}\right)$, we can define an automorphism $\sigma^{\prime}$ of $\mathfrak{K}\left[z_{2}\right]=\mathfrak{K}\left[z_{0}\right]$ such that $t\left(a, z_{2}\right)^{\sigma^{\prime}}=t\left(a y^{-1}, z_{2}\right)$; then $\mathcal{P}\left(z_{2}\right)^{\sigma^{\prime}}=\mathcal{P}\left(z_{2}\right)$ by Lemma 7.8. Put $\varepsilon^{\prime}=\varepsilon^{\sigma^{\prime}}$. Then

$$
\begin{equation*}
t\left(a, z_{1}\right)^{\sigma \sigma^{\prime}}=\left(\varepsilon t\left(a y, z_{2}\right)\right)^{\sigma^{\prime}}=\varepsilon^{\prime} t\left(a, z_{2}\right) \tag{**}
\end{equation*}
$$

Now let $N^{\prime}$ be a multiple of $N$; taking $N^{\prime}$ to be $N$ of $\S 7.9$, consider $g=$ $\left(g_{1}, \ldots, g_{m}\right), M, V$, and $q$ as in $\S 7.9$. We can find a finitely generated extension $k^{\prime}$ of $\mathbf{Q}$ contained in $k$ such that $V_{M}, V_{N^{\prime}}, p_{N^{\prime}}^{\Lambda \prime}, V$, and $q$ are all $k^{\prime}$-rational and $g_{i} \in k^{\prime}\left(V_{M}\right)$. Since $\mathfrak{K}=\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, by Lemma 7.4 (1), $\mathfrak{K}$ and $k^{\prime} \mathbf{Q}_{\mathrm{ab}}$ are linearly disjoint over $\mathbf{Q}_{\mathrm{ab}}$, and hence $\mathfrak{K}\left[z_{0}\right]$ and $k^{\prime} \mathbf{Q}_{\mathrm{ab}}$ are linearly disjoint over $\mathbf{Q}_{\mathrm{ab}}$.

Therefore the automorphism $\sigma \sigma^{\prime}$ of $\mathfrak{K}\left[z_{0}\right]$ over $\mathbf{Q}_{\mathrm{ab}}$ can be extended to an automorphism $\tau$ of $k^{\prime} \mathfrak{K}\left[z_{0}\right]$ over $k^{\prime} \mathbf{Q}_{\mathrm{ab}}$. Clearly $z_{\nu}$ is generic for $k^{\prime} \mathfrak{F}$ over $k^{\prime}$. Now $\varphi_{M}\left(z_{1}\right)^{\tau}=\varphi_{M}(w)$ with $w \in \mathcal{H}$ and $g\left(z_{1}\right)^{\tau}=q\left(\varphi_{M}\left(z_{1}\right)\right)^{\tau}=q\left(\varphi_{M}(w)\right)=g(w)$. From this we can derive that $\left(\mathcal{P}_{z_{1}}^{N^{\prime}}\right)^{\tau}=\mathcal{P}_{w}^{N^{\prime}}$. Indeed, $w$ is generic for $g$ over $k^{\prime}$ and so generic for $\mathfrak{K}$ over $k^{\prime}$, since $\mathfrak{K}$ is algebraic over $\mathbf{Q}(g)$. Thus $f\left(z_{1}\right) \mapsto f(w)$ for $f \in \mathfrak{K}$ is an isomorphism, which coincides with $\tau$ on $\mathbf{Q}\left(g\left(z_{1}\right)\right)$, and which sends $\mathcal{P}_{z_{1}}^{N^{\prime}}$ to $\mathcal{P}_{w}^{N^{\prime}}$ by Lemma 7.8. Since $\mathcal{P}_{z_{1}}^{N^{\prime}}$ is defined over $\mathbf{Q}\left(g\left(z_{1}\right)\right)$, this shows that $\left(\mathcal{P}_{z_{1}}^{N^{\prime}}\right)^{\tau} .=\mathcal{P}_{w}^{N^{\prime}}$, so that $t\left(a, z_{1}\right)^{\tau}=t(a, w)$ for every $a \in\left(N^{\prime}\right)^{-1} L$. Since $\mathcal{P}\left(z_{2}\right)^{\sigma^{\prime}}=\mathcal{P}\left(z_{2}\right),\left({ }^{* *}\right)$ together with (7.5a) shows that $\varepsilon^{\prime}$ is an isomorphism of $\mathcal{P}_{z_{2}}^{N^{\prime}}$ to $\left(\mathcal{P}_{z_{1}}^{N^{\prime}}\right)^{\sigma \sigma^{\prime}}$. Thus $z_{2}=\gamma w$ with some $\gamma \in \Gamma^{N^{\prime}}$, and for $a \in\left(N^{\prime}\right)^{-1} L$ we have $t\left(a, z_{2}\right)=t(a, \gamma w)=t(a \gamma, w)=t(a, w)=t\left(a, z_{1}\right)^{\tau}$. This shows that $t\left(a, z_{1}\right)^{\sigma \sigma^{\prime}}=t\left(a, z_{2}\right)$ for every $a \in K_{2 n}^{1}$, since our reasoning is applicable to every multiple $N^{\prime}$ of $N$. Thus $t\left(a, \alpha z_{0}\right)^{\sigma \sigma^{\prime}}=t\left(a, z_{1}\right)^{\sigma \sigma^{\prime}}=t\left(a, z_{2}\right)=t\left(a y, z_{2}\right)^{\sigma^{\prime}}$. Applying $\sigma^{\prime-1}$ to this, we obtain the expected equality.

It remains to prove that $\varepsilon$ is an isomorphism of $\mathcal{P}\left(z_{2}\right)$ onto $\mathcal{P}\left(z_{1}\right)^{\sigma}$. We first observe that $\nu\left(y x^{-1}\right)=\nu\left(\alpha_{1} \alpha_{2}^{-1}\right) \in F \cap F_{\mathrm{a}+} \mathbf{Z}_{\mathrm{h}}^{\times}$, so that $\nu(x)=\nu(y)$. Applying $\sigma$ to (7.5a) and employing (*), we obtain $\iota_{z_{1}}(c)^{\sigma} \varepsilon t\left(a y, z_{2}\right)=\iota_{z_{1}}(c)^{\sigma} t\left(a, z_{1}\right)^{\sigma}=$ $t\left(c a, z_{1}\right)^{\sigma}=\varepsilon t\left(c a y, z_{2}\right)=\varepsilon \iota_{z_{2}}(c) t\left(a y, z_{2}\right)$, so that $\iota_{z_{1}}(c)^{\sigma} \varepsilon=\varepsilon \iota_{z_{2}}(c)$. Next, let $X_{\nu}$ denote the divisor $X_{z}$ for $z=z_{\nu}$ given in the proof of Lemma 7.12. By (7.10) we have $\zeta_{X_{\nu}}\left(t\left(a, z_{\nu}\right), t\left(b, z_{\nu}\right)\right)=\mathbf{e}\left((m / s) \operatorname{Tr}_{K / \mathbf{Q}}\left(a \eta_{n} b^{*}\right)\right)$ for every $a, b \in m^{-1} L, 0<$ $m \in \mathbf{Z}$. Let $Y=\varepsilon^{-1}\left(X_{1}^{\sigma}\right)$ and let $E^{\prime}$ be the Riemann form determined by $Y$. Then by $\left(^{*}\right),(2.10)$, and (2.11) we have, for such $a$ and $b$,

$$
\begin{gathered}
\mathbf{e}\left\{m E^{\prime}\left(p_{z_{2}}(a y), p_{z_{2}}(b y)\right)\right\}=\zeta_{Y}\left(t\left(a y, z_{2}\right), t\left(b y, z_{2}\right)\right)=\zeta_{X_{1}^{\sigma}}\left(\varepsilon t\left(a y, z_{2}\right), \varepsilon t\left(b y, z_{2}\right)\right) \\
=\zeta_{X_{1}}\left(t\left(a, z_{1}\right), t\left(b, z_{1}\right)\right)^{\sigma}=\mathbf{e}\left((m / s) \operatorname{Tr}_{K / \mathbf{Q}}\left(a \eta_{n} b^{*}\right)\right)^{\sigma}
\end{gathered}
$$

In view of (8.2) and (7.10), the last quantity equals $\mathbf{e}\left\{m E_{z_{2}}\left(p_{z_{2}}(a y), p_{z_{2}}(b y)\right)\right\}$, since $\sigma=\left[\nu(x)^{-1}, \mathbf{Q}\right]$ on $\mathbf{Q}_{\mathrm{ab}}$ and $\nu(x)=\nu(y)$. Thus $E^{\prime}=E_{z_{2}}$, which means that $\varepsilon$ sends $\mathcal{C}_{z_{2}}$ to $\left(\mathcal{C}_{z_{1}}\right)^{\sigma}$. This completes the proof.
8.8. Put

$$
\begin{equation*}
\mathcal{G}=\left(G_{0}\right)_{\mathbf{A}} \widetilde{G} \widetilde{G}_{\mathbf{a}+}, \quad \mathcal{G}_{+}=\mathcal{G} \cap \widetilde{G}_{\mathbf{A}+} \tag{8.9}
\end{equation*}
$$

Since $\left(G_{0}\right)_{\mathbf{A}}$ and $\widetilde{G}_{\mathbf{a}+}$ are normal subgroups of $\widetilde{G}_{\mathbf{A}}, \mathcal{G}$ and $\mathcal{G}_{+}$are subgroups of $\widetilde{G}_{\mathrm{A}} ;$ moreover,

$$
\begin{equation*}
\mathcal{G}_{+}=\left(G_{0}\right)_{\mathbf{A}+} \widetilde{G}_{+} \widetilde{G}_{\mathbf{a}+}=T^{N} \widetilde{G}_{+} \widetilde{G}_{\mathbf{a}+}=\widetilde{G}_{+} T^{N} \widetilde{G}_{\mathbf{a}+} \text { for every } \quad N \tag{8.10}
\end{equation*}
$$

Indeed, let $\alpha \in\left(G_{0}\right)_{\mathbf{A}}$ and $\beta \in \widetilde{G}$; suppose $\nu(\alpha \beta) \gg 0$. Since $\nu(\alpha) \in \mathbf{Q}_{\mathbf{A}}^{\times}$, we can find an element $\gamma$ of $G_{0}$ such that that $\nu(\alpha \gamma) \gg 0$. Then $\alpha \beta=\alpha \gamma \gamma^{-1} \beta$, $\alpha \gamma \in\left(G_{0}\right)_{\mathbf{A}+}$, and $\gamma^{-1} \beta \in \widetilde{G}_{+}$. The equalities of (8.10) follow from this fact and Lemma 8.3 (1).

Now given $x \in \mathcal{G}_{+}$, take $c \in \mathbf{Q}_{\mathbf{A}}^{\times}$so that $\nu(x) \in c F^{\times} F_{\mathbf{a}+}^{\times}$. We then define $\sigma(x)$ to be $\left[c^{-1}, \mathbf{Q}\right]\left(\in \operatorname{Aut}\left(\mathbf{Q}_{\mathrm{ab}}\right)\right)$. Then we easily see that $\sigma(x)$ is well-defined independently of the choice of $c$, and thus we obtain a homomorphism

$$
\begin{equation*}
\sigma: \mathcal{G}_{+} \longrightarrow \operatorname{Gal}\left(\mathbf{Q}_{\mathrm{ab}} / \mathbf{Q}\right) \tag{8.11}
\end{equation*}
$$

8.9. Lemma. (1) $\alpha U^{N} \alpha^{-1}=U^{N}$ for every $N$ if $\alpha \in \widetilde{G}_{\mathbf{A}}$ and $L \alpha=L$.
(2) $\Gamma^{N} U^{N}$ and $\Gamma^{N} T^{N}$ are subgroups of $\widetilde{G}_{\mathbf{A}}$.
(3) $\Gamma^{N} U^{N}$ is a normal subgroup of $\Gamma^{1} T^{1}$.
(4) $\Gamma^{N} T^{N} / \Gamma^{N} U^{N}$ is isomorphic to $T^{N} / U^{N}$.

Proof. Since $\left(G_{0}\right)_{\mathbf{A}}$ is normal in $\widetilde{G}_{\mathbf{A}}$, we easily obtain (1). In particular $\gamma U^{N} \gamma^{-1}=U^{N}$ for every $\gamma \in \Gamma^{N}$, so that $\Gamma^{N} U^{N}$ is a subgroup of $\widetilde{G}_{\mathbf{A}}$. Let $x \in T^{1}$ and $\gamma \in \Gamma^{N}$. Then $x^{-1} \gamma x \gamma^{-1}-1 \prec N \mathfrak{r}$, so that $x^{-1} \gamma x \gamma^{-1} \in U^{N}$, and hence $\gamma x \gamma^{-1} \in T^{N}$ if $x \in T^{N}$. Thus $\Gamma^{N} T^{N}$ is a subgroup of $\widetilde{G}_{\mathbf{A}}$. Clearly $\alpha\left(\Gamma^{N} U^{N}\right) \alpha^{-1}=\Gamma^{N} U^{N}$ for $\alpha \in \Gamma^{1}$. Now let $x \in T^{1}$. Then $x U^{N} x^{-1}=U^{N}$. For $\gamma \in \Gamma^{N}$, we have $x^{-1} \gamma x \in U^{N} \gamma \subset \Gamma^{N} U^{N}$. Thus $x \Gamma^{N} x^{-1} \subset \Gamma^{N} U^{N}$. This proves (3). From Lemma 8.3 (2) we obtain $\Gamma^{N} U^{N} \cap T^{N}=U^{N}$, from which (4) follows.
8.10. Theorem. There exists a homomorphism $\tau: \mathcal{G}_{+} \longrightarrow \operatorname{Aut}(\mathfrak{K})$ with the following properties:
(1) $\tau(\xi)=\sigma(\xi)$ on $\mathbf{Q}_{\mathrm{ab}}$.
(2) $f^{\tau(\alpha)}=f \circ \alpha$ for every $f \in \mathfrak{K}$ and every $\alpha \in \widetilde{G}_{+}$.
(3) $f^{\tau(\xi)}=f^{[s, \mathbf{Q}]}$ if $\xi=\iota(s)$ with $s \in \mathbf{Z}_{\mathbf{h}}^{\times}$.
(4) $\operatorname{Ker}(\tau)=K^{\times} \widetilde{G}_{\mathbf{a}+}$.
(5) For any fixed $f \in \mathfrak{K}$ the set $\left\{\xi \in \mathcal{G}_{+} \mid f^{\tau(\xi)}=f\right\}$ contains $U^{N}$ for some $N$.
(6) Let $k_{N}=\mathbf{Q}(\mathbf{e}(1 / N))$. Then, for every $N$ we have

$$
\begin{align*}
\mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right) & =\left\{f \in \mathfrak{K} \mid f^{\tau(x)}=f \quad \text { for every } x \in \Gamma^{N} U^{N}\right\}  \tag{8.12}\\
\mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right) & =\left\{f \in \mathfrak{K} \mid f^{\tau(x)}=f \quad \text { for every } x \in \Gamma^{N} T^{N}\right\} \tag{8.13}
\end{align*}
$$

Here we remind the reader that $\mathfrak{K}=\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$.
Proof. Fix any $N$ as in (8.8). By (8.10), given $\xi \in \mathcal{G}_{+}$, we can find $x \in T^{N}$ and $\alpha \in \widetilde{G}_{+}$such that $\xi \in x \alpha \widetilde{G}_{\mathbf{a}+}$ with $x \in T^{N}$ and $\alpha \in \widetilde{G}_{+}$. We then define an automorphism $\underset{\sim}{\tau}(\xi)$ of $\mathfrak{K}$ by $f^{\tau(\xi)}=f^{\tau(x)} \circ \alpha$. To show that this is well-defined, let $x \alpha \in x_{1} \alpha_{1} \widetilde{G}_{\mathbf{a}+}$ with $x_{1} \in T^{M}$ and $\alpha_{1} \in \widetilde{G}_{+}$with a multiple $M$ of $N$. Put $\gamma=\alpha_{1} \alpha^{-1}$. By Lemma $8.3(2), \gamma \in \Gamma^{N} \cap T^{N}$, and $f^{\tau(x)} \circ \alpha=f^{\tau\left(x_{1} \gamma\right)} \circ \alpha=$ $f^{\tau\left(x_{1}\right)} \circ \gamma \alpha=f^{\tau\left(x_{1}\right)} \circ \alpha_{1}$, since $g^{\tau(\gamma)}=g \circ \gamma$ as shown in Lemma 8.6. Thus $\tau(\xi)$ is well-defined independently of $(N, x, \alpha)$. Next let $\xi^{\prime} \in x^{\prime} \alpha^{\prime} \widetilde{G}_{\mathrm{a}+}$ with $x^{\prime} \in T^{N}$ and $\alpha^{\prime} \in \widetilde{G}_{+}$. By (8.10) we have $\alpha x^{\prime} \in y \beta \widetilde{G}_{\mathbf{a}+}$ with $y \in T^{N}$ and $\beta \in \widetilde{G}_{+}$. Then $\xi \xi^{\prime} \in x y \beta \alpha^{\prime} \widetilde{G}_{\mathbf{a}+}$, and employing Lemma 8.7 , we can easily verify that $\tau(\xi) \tau\left(\xi^{\prime}\right)=$ $\tau\left(\xi \xi^{\prime}\right)$. Now $\tau(\xi)=\tau(x)$ on $\mathbf{Q}_{\mathrm{ab}}$, and $\tau(x)=\left[\nu(x)^{-1}, \mathbf{Q}\right]$ on $\mathbf{Q}_{\mathrm{ab}}$ by Lemma 8.6. Since $\nu(\xi) \in \nu(x) F^{\times} F_{\mathbf{a}+}^{\times}$, this proves (1). Property (3) is clear from our definition and Lemmas $8.5,8.6$. To prove (4), let $\xi \in \alpha \widetilde{G}_{\mathbf{a}+}$ with $\alpha \in K^{\times}$. Then $t(a, z)^{\tau(\xi)}=t(a, \alpha z)=t(a, z)$, so that $\tau(\xi)=$ id. Suppose conversely $\tau(\xi)=$ id. on $\mathfrak{K}$ and $\xi \in x \alpha \widetilde{G}_{\mathbf{a}+}$ with $x \in T^{N}$ and $\alpha \in \widetilde{G}_{+}$. Then

$$
\begin{equation*}
t(a, z)=t(a, z)^{\tau(\xi)}=t(a x, \alpha z) \quad \text { for every } \quad a \in K_{2 n}^{1} \tag{8.14}
\end{equation*}
$$

Since $\sigma(\xi)=1$, we have $\left[\nu(x)^{-1}, \mathbf{Q}\right]=\sigma(x)=1$. Now $\nu(x) \in \nu\left(T^{N}\right)=\mathbf{Z}_{\mathbf{h}}^{\times} \mathbf{Q}_{\mathbf{a}+}^{\times}$, so that $\nu(x)_{\mathbf{h}}=1$. Then we easily see that $x \in U^{N}$. Fix a generic point $z_{0} \in \mathcal{H}$ as before. Then (8.14) shows that $t\left(a, z_{0}\right)=t\left(a x, \alpha z_{0}\right)$, and hence $\mathcal{P}^{N}\left(z_{0}\right)=$ $\mathcal{P}^{N}\left(\alpha z_{0}\right)$ by Lemma 7.8 , since $x-1 \prec N \mathfrak{r}$. Thus $\alpha z_{0}=\gamma z_{0}$ with $\gamma \in \Gamma^{N}$. Taking $z_{0}$ generic even for the action of $\widetilde{G}_{+}$, we find that $\alpha \in K^{\times} \gamma$. Now $t(a, z)=$ $t(a x, \gamma z)=t(a x \gamma, z)$ for every $a \in K_{2 n}^{1}$, and hence $x \gamma \in \widetilde{G}_{\mathbf{a}+}$, so that $\xi=x \alpha \in$ $K^{\times} \widetilde{G}_{\mathbf{a}+}$. This proves (4). To prove (5), we observe that $t(a, z)^{\tau(x)}=t(a x, z)=$
$t(a, z)$ if $x \in U^{N}$ and $a \in N^{-1} L$. Since $\mathfrak{K}$ is generated by the coordinates of such points $t(a, z)$, we easily obtain (5).

To prove (6) (in which we do not assume (8.8)), let us write simply $f^{x}$ for $f^{\tau(x)}$. Let $f$ be an element of the right-hand side of (8.12). Then $f \circ \gamma=f$ for every $\gamma \in \Gamma^{N}$. Moreover, if $s \in \mathbf{Z}_{\mathbf{h}}^{\times}$and $s-1 \prec N \mathbf{Z}$, then $\iota(s) \in U^{N}$, and so $f^{\iota(s)}=f$. Thus, by Lemma $8.5(3), f \in \mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$. Similarly if $f$ belongs to the right-hand side of (8.13), then the same type of argument shows that $f \in \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$. Conversely, let $f \in \mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$. By (5) there exists a multiple $M$ of $N$ such that $f^{\xi}=f$ for every $\xi \in U^{M}$. Given $x \in U^{N}$, by Lemma 8.3 (4) we can put $x=\iota(r) y \gamma$ with $r=\nu(x)^{-1}, y \in U^{M}$, and $\gamma \in \Gamma^{N}$. Since $f \in \mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$, we have $f^{\iota(r)}=f=f \circ \gamma$, and hence $f^{x}=f^{\iota(r) y \gamma}=f^{y} \circ \gamma=f$, which proves (8.12). Next, let $f \in \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$ and $x \in T^{N}$. Then $x=\iota(p) x^{\prime}$ with $p \in \mathbf{Z}_{\mathbf{h}}^{\times}$and $x^{\prime} \in U^{N}$. Since $f^{\iota(p)}=f$ and $f^{x^{\prime}}=f$ by (8.12), we obtain $f^{x}=f$. This proves (8.13).
8.11. Theorem. (1) $\mathfrak{K}$ is algebraic over $\mathcal{A}_{0}\left(\Gamma^{1}, \mathbf{Q}\right)$.
(2) $\mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$ and $\mathbf{C}$ are linearly disjoint over $\mathbf{Q}$. Moreover, $\mathcal{A}_{0}\left(\Gamma^{N}, k\right)=$ $k \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$ for every subfield $k$ of $\mathbf{C}$; in particular $\mathcal{A}_{0}\left(\Gamma^{N}\right)=\mathbf{C} \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$.
(3) $\mathcal{A}_{0}(\mathbf{Q})$ and $\mathbf{C}$ are linearly disjoint over $\mathbf{Q}$. Moreover $\mathcal{A}_{0}(k)=k \mathcal{A}_{0}(\mathbf{Q})$ for every subfield $k$ of $\mathbf{C}$; in particular $\mathcal{A}_{0}=\mathbf{C} \mathcal{A}_{0}(\mathbf{Q})$.
(4) Given $f \in \mathcal{A}_{0}$, there exists a finitely generated extension $k$ of $\mathbf{Q}$ such that $f \in \mathcal{A}_{0}(k)$.
(5) Given $f \in \mathcal{A}_{0}$ and $\sigma \in \operatorname{Aut}(\mathbf{C})$, the element $f^{\sigma}$ formally defined in $\S 5.9$ is indeed an element of $\mathcal{A}_{0}$.

Proof. The linear disjointness in (2) and (3) follows immediately from Lemma 7.4 (1). Now observe that $T^{N} / U^{N}$ is isomorphic to $(\mathbf{Z} / N \mathbf{Z})^{\times}$(resp. to $\operatorname{Gal}\left(k_{N} / \mathbf{Q}\right)$ ) via the map $x \mapsto \nu(x)$ (resp. via $x \mapsto \sigma(x))$ for $x \in T^{N}$. From Lemma 8.9 (3) and (8.12) we see that $\mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$ is stable under $\tau\left(T^{N}\right)$. Thus we obtain a homomorphism of $T^{N} / U^{N}$ into the group of automorphisms of $\mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$. This is injective, since $T^{N} / U^{N}$ is isomorphic to $\operatorname{Gal}\left(k_{N} / \mathbf{Q}\right)$. By (8.13) the fixed subfield of $T^{N} / U^{N}$ is $\mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$, so that $\mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$ is a Galois extension of $\mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$ whose Galois group is isomorphic to $T^{N} / U^{N}$. By the linear disjointness we have $\left[k_{N} \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right): \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)\right]=\left[k_{N}: \mathbf{Q}\right]$, and hence we obtain

$$
\begin{equation*}
\mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)=k_{N} \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right) \tag{8.15}
\end{equation*}
$$

Now $\mathcal{A}_{0}=\mathbf{C} \mathfrak{K}$ by Theorem 7.10 (3) and every quotient of (7.8) belongs to $\mathcal{A}_{0}\left(\Gamma^{N}\right.$, $k_{N}$ ) for some $N$. Therefore we can conclude that $\mathcal{A}_{0}=\mathbf{C} \mathcal{A}_{0}(\mathbf{Q})$. Then (4) and (5) follow immediately from this fact. Now, for any subfield $k$ of $\mathbf{C}$, we have $\mathbf{C} \mathcal{A}_{0}(k)=\mathcal{A}_{0}=\mathbf{C} k \mathcal{A}_{0}(\mathbf{Q})$, and $k \mathcal{A}_{0}(\mathbf{Q}) \subset \mathcal{A}_{0}(k)$. Since $\mathbf{C}$ and $\mathcal{A}_{0}(k)$ are linearly disjoint over $k$, we obtain $\mathcal{A}_{0}(k)=k \mathcal{A}_{0}(\mathbf{Q})$ by Lemma 7.2. This completes the proof of (3).

To prove the main part of (2), given $N$, take a finitely generated extension $k$ of Q so that $V_{N}$ is defined over $k$. Then $k\left(V_{N}\right) \circ \varphi_{N}=k\left(f_{1}, \ldots, f_{m}\right)$ with suitable $f_{i} \in \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{C}\right)$. By (4) we can find a finitely generated extension $k^{\prime}$ of $k$ so that $f_{i} \in \mathcal{A}_{0}\left(\Gamma^{N}, k^{\prime}\right)$ for every $i$. Then $k^{\prime}\left(V_{N}\right) \circ \varphi_{N}=k^{\prime}\left(f_{1}, \ldots, f_{m}\right) \subset \mathcal{A}_{0}\left(\Gamma^{N}, k^{\prime}\right)$. By Lemma $7.4(2)$ we have $k^{\prime}\left(f_{1}, \ldots, f_{m}\right)=\mathcal{A}_{0}\left(\Gamma^{N}, k^{\prime}\right)$. Since $\mathcal{A}_{0}\left(k^{\prime}\right)=k^{\prime} \mathcal{A}_{0}(\mathbf{Q})$, we can find $g_{1}, \ldots, g_{r} \in \mathcal{A}_{0}(\mathbf{Q})$ so that $f_{1}, \ldots, f_{m} \in k^{\prime}\left(g_{1}, \ldots, g_{r}\right)$. Take a multiple $M$ of $N$ so that $g_{1}, \ldots, g_{r} \in \mathcal{A}_{0}\left(\Gamma^{M}, \mathbf{Q}\right)$. Replacing $k^{\prime}$ by $k^{\prime} k_{M}$, we may assume that $k_{M} \subset k^{\prime}$. From (8.12) and Lemma 8.9 (3) we see that $\mathcal{A}_{0}\left(\Gamma^{M}, k_{M}\right)$
is stable under $\Gamma^{N}$. Since $\mathcal{A}_{0}\left(\Gamma^{N}, k_{M}\right)$ consists of the $\Gamma^{N}$-invariant elements of $\mathcal{A}_{0}\left(\Gamma^{M}, k_{M}\right)$, we see that $\mathcal{A}_{0}\left(\Gamma^{M}, k_{M}\right)$ is a Galois extension of $\mathcal{A}_{0}\left(\Gamma^{N}, k_{M}\right)$. By Lemma $7.4(1), \mathcal{A}_{0}\left(k_{M}\right)$ and $k^{\prime}$ are linearly disjoint over $k_{M}$, so that $k^{\prime} \mathcal{A}_{0}\left(\Gamma^{M}, k_{M}\right)$ is a Galois extension of $k^{\prime} \mathcal{A}_{0}\left(\Gamma^{N}, k_{M}\right)$ with the same Galois group. Now the elements of $\mathcal{A}_{0}\left(\Gamma^{N}, k^{\prime}\right)$ are contained in $k^{\prime} \mathcal{A}_{0}\left(\Gamma^{M}, k_{M}\right)$ and $\Gamma^{N_{\text {N }} \text {-invariant, so that }}$ $\mathcal{A}_{0}\left(\Gamma^{N}, k^{\prime}\right)=k^{\prime} \mathcal{A}_{0}\left(\Gamma^{N}, k_{M}\right)$. Now we can let $\Gamma^{N} U^{N} / \Gamma^{M} U^{M}$ act on $\mathcal{A}_{0}\left(\Gamma^{M}, k_{M}\right)$, and observe that $\Gamma^{N} U^{N} / \Gamma^{M} U^{M}$ can be mapped onto the Galois group of $\mathcal{A}_{0}\left(\Gamma^{M}\right.$, $k_{M}$ ) over $\mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$. (Taking $N=1$ here, we obtain (1).) We easily see that $\tau(\alpha x)$, with $\alpha \in \Gamma^{N}$ and $x \in U^{N}$, gives the identity map on $k_{M} \mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$ if and only if $\nu(x)-1 \prec M \mathbf{Z}$. Given such an $x$, by Lemma 8.3 (4) we have $x=\iota(r) y \gamma$ with $r=\nu(x)^{-1}, y \in U^{M}$, and $\gamma \in \Gamma^{N}$. Then $[r, \mathbf{Q}]=$ id. on $k_{M}$, so that for $f \in \mathcal{A}_{0}\left(\Gamma^{N}, k_{M}\right)$ and $\alpha \in \Gamma^{N}$ we have $f^{\tau(\alpha x)}=f^{\tau(y \gamma)}=f$. By Galois theory this shows that $\mathcal{A}_{0}\left(\Gamma^{N}, k_{M}\right)=k_{M} \mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$. This combined with (8.15) shows that $\mathcal{A}_{0}\left(\Gamma^{N}, k_{M}\right)=k_{M} \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$. Thus $\mathcal{A}_{0}\left(\Gamma^{N}, k^{\prime}\right)=$ $k^{\prime} \mathcal{A}_{0}\left(\Gamma^{N}, k_{M}\right)=k^{\prime} \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$. Consequently $\mathcal{A}_{0}\left(\Gamma^{N}\right)=\mathbf{C}\left(V_{N}\right) \circ \varphi_{N}=\mathbf{C} k^{\prime}\left(V_{N}\right) \circ$ $\varphi_{N}=\mathbf{C} \mathcal{A}_{0}\left(\Gamma^{N}, k^{\prime}\right)=\mathbf{C} \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$. This proves the last statement of (2). Then for any subfield $k_{0}$ of $\mathbf{C}$ we have $\mathbf{C} \mathcal{A}_{0}\left(\Gamma^{N}, k_{0}\right) \subset \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{C}\right)=\mathbf{C} \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)=$ $\mathbf{C} k_{0} \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right) \subset \mathbf{C} \mathcal{A}_{0}\left(\Gamma^{N}, k_{0}\right)$, so that $\mathbf{C} k_{0} \cdot \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)=\mathbf{C} \mathcal{A}_{0}\left(\Gamma^{N}, k_{0}\right)$. Since $\mathcal{A}_{0}\left(\Gamma^{N}, k_{0}\right)$ and $\mathbf{C}$ are linearly disjoint over $k_{0}$ by Lemma 7.4 (1), we obtain $k_{0} \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)=\mathcal{A}_{0}\left(\Gamma^{N}, k_{0}\right)$ by Lemma 7.2. This completes the proof.
8.12. Theorem. For a point $w$ on $\mathcal{H}$ and a positive integer $N$, let $\mathfrak{\kappa}_{N}[w]$ denote the field consisting of $f(w)$ for all $f \in \mathfrak{K}_{N}$, where $\mathfrak{K}_{N}=\mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$. Then $\mathfrak{K}_{N}[w]$ is the field of moduli of $\mathcal{P}_{w}^{N}$.

Proof. We shall prove this in the proof of Theorem 9.3, invoking some results of [S66a]. Here we prove the case where $w$ is generic for $\mathfrak{K}$ over $\mathbf{Q}$, using the results obtained in this section so far. In $\S 7.9$ we observed that $\mathcal{P}_{w}^{N}$ is rational over $\mathfrak{K}[w]$. Therefore we can find a multiple $M$ of $N$ such that $\mathcal{P}_{w}^{N}$ is rational over $\mathfrak{K}_{M}[w]$. Since $f \mapsto f(w)$ is an isomorphism of $\mathfrak{K}$ onto $\mathfrak{K}[w]$, we can let $\mathcal{G}_{+}$act on $\mathfrak{K}[w]$ by putting $f(w)^{x}=f^{\tau(x)}(w)$ for $x \in \mathcal{G}_{+}$and $f \in \mathfrak{K}$. As shown in the proof of Theorem 8.11, $\Gamma^{N} U^{N} / \Gamma^{M} U^{M}$ is mapped onto $\operatorname{Gal}\left(\mathfrak{K}_{M}[w] / \mathfrak{K}_{N}[w]\right)$. Let $\Phi$ be the field of moduli of $\mathcal{P}_{w}^{N}$. To show that $\Phi \subset \mathfrak{K}_{N}[w]$, take $x \in U^{N} \Gamma^{N}$. By Lemma 8.3 (4), $x=\iota(r) y \gamma$ with $r=\nu(x)^{-1}, y \in U^{M}$, and $\gamma \in \Gamma^{N}$. Since $t(a, w)^{\iota(r)}=t(a \iota(r), w)$ and $r-1 \prec N \mathbf{Z}$, we obtain $\left(\mathcal{P}_{w}^{N}\right)^{\iota(r)}=\mathcal{P}_{w}^{N}$ by Lemma 7.8. Since $\mathcal{P}_{w}^{N}$ is rational over the field $\mathfrak{K}_{M}[w]$ on which $y$ gives the identity map, we have $\left(\mathcal{P}_{w}^{N}\right)^{y}=\mathcal{P}_{w}^{N}$. Thus $\left(\mathcal{P}_{w}^{N}\right)^{x}=\left(\mathcal{P}_{w}^{N}\right)^{\gamma}=\mathcal{P}_{\gamma w}^{N}$, as $f(w)^{\gamma}=f(\gamma w)$ for every $f \in \mathfrak{K}$. Since $\gamma \in \Gamma^{N}$, $\mathcal{P}_{\gamma w}^{N}$ is isomorphic to $\mathcal{P}_{w}^{N}$. Therefore the property of $\Phi$ stated in Theorem 2.8 (1), (ii) implies that $x$ gives the identity map on $\Phi$, so that $\Phi \subset \mathfrak{K}_{N}[w]$.

To show that $\mathfrak{K}_{N}[w] \subset \Phi$, we consider $V_{N}, \varphi_{N}$, and $p_{N}^{M}$ as in $\S 7.9$. We take a finitely generated extension $k$ of $k_{M}$ over which $V_{M}, V_{N}$, and $p_{N}^{M}$ are rational. Since $k\left(V_{M}\right) \circ \varphi_{M} \subset \mathcal{A}_{0}\left(\Gamma^{M}\right) \subset \mathbf{C} \mathcal{A}_{0}\left(\Gamma^{M}, \mathbf{Q}\right)$, changing $k$ suitably, we may assume that $k\left(V_{M}\right) \circ \varphi_{M} \subset \mathcal{A}_{0}\left(\Gamma^{M}, k\right)$. Then $k\left(V_{M}\right) \circ \varphi_{M}=\mathcal{A}_{0}\left(\Gamma^{M}, k\right)$ by Lemma 7.4 (2). Since $\varphi_{N}=p_{N}^{M} \circ \varphi_{M}$, we obtain $k\left(V_{N}\right) \circ \varphi_{N}=\mathcal{A}_{0}\left(\Gamma^{N}, k\right)$. Now $k_{N} \subset \Phi$ by Lemma 7.12. To prove our theorem, we may assume that $w$ is generic for $k \mathfrak{K}$ over $k$. Then $\varphi_{N}(w)$ (resp. $\varphi_{M}(w)$ ) is a generic point of $V_{N}$ (resp. $V_{M I}$ ) over $k$. Let $\sigma$ be an isomorphism of $\mathfrak{K}_{N}[w]$ onto a subfield of $\mathbf{C}$ over $\Phi$. Since $k$ and $\mathfrak{K}_{N}[w]$ are linearly disjoint over $k_{N}$, we can extend $\sigma$ to an isomorphism of $k \mathfrak{K}_{N}[w]$ onto a subfield of $\mathbf{C}$ over $k \Phi$. Extend this further to $k \mathfrak{K}_{M I}[w]$. Then $\varphi_{M I}(w)^{\sigma}$ is a point
of $V_{M}$, so that $\varphi_{M}(w)^{\sigma}=\varphi_{M}\left(w^{\prime}\right)$ with some $w^{\prime} \in \mathcal{H}$. Since $p_{N}^{M}$ is $k$-rational, we have $\varphi_{N}(w)^{\sigma}=\varphi_{N}\left(w^{\prime}\right)$, and $w^{\prime}$ is generic for $k \mathfrak{K}_{M}$ over $k$. Since $\mathcal{P}_{w}^{N}$ is rational over $k \mathfrak{K}_{M}[w]=k\left(\varphi_{M}(w)\right)$, we have $\left(\mathcal{P}_{w}^{N}\right)^{\sigma}=\mathcal{P}_{w^{\prime}}^{N}$. Since $\sigma=$ id. on $\Phi,\left(\mathcal{P}_{w}^{N}\right)^{\sigma}$ must be isomorphic to $\mathcal{P}_{w}^{N}$, so that $w^{\prime} \in \Gamma^{N} w$. Therefore $\varphi_{N}(w)^{\sigma}=\varphi_{N}\left(w^{\prime}\right)=\varphi_{N}(w)$, and hence $\sigma=$ id. on $\mathfrak{K}_{N}[w]$. Thus $\mathfrak{K}_{N}[w] \subset \Phi$, which completes the proof.

## 9. The reciprocity-law at CM-points and rationality of automorphic forms

The first main purpose of this section is to study the behavior of the values of the elements of $\mathfrak{K}$ at a CM-point under certain automorphisms, which may be viewed as an explicit reciprocity law of a certain abelian extension. The next task is to extend the results of the previous two sections to automorphic forms of nontrivial weight. First we quote a theorem without proof:
9.1. Theorem. Given a PEL-type $\Omega=\left\{K, \Psi, L, \mathcal{T},\left\{u_{i}\right\}_{i=1}^{s}\right\}$ as in (4.7) (in all three cases), consider the family $\mathcal{F}(\Omega)=\left\{\mathcal{P}_{z} \mid z \in \mathcal{H}\right\}$ as in (4.26) and define $\Gamma$ by (4.28). Then there exist an algebraic number field $k_{\Omega}$ of finite degree and a model $(V, \varphi)$ of $\Gamma \backslash \mathcal{H}$ with the following properties:
(1) Let $K_{\Psi}$ be the field generated over $\mathbf{Q}$ by the numbers $\operatorname{tr}(\Psi(c))$ for all $c \in K$. Then $K_{\Psi} \subset k_{\Omega}$.
(2) If $\alpha \in \operatorname{Aut}(\mathbf{C})$ and $\mathcal{P}$ is a structure of type $\Omega$, then $\mathcal{P}^{\alpha}$ is of type $\Omega$ if and only if $\alpha=$ id. on $k_{\Omega}$, where $\mathcal{P}^{\alpha}$ is defined as in $\S 2.7$.
(3) $V$ is defined over $k_{\Omega}$.
(4) $k_{\Omega}(\varphi(w))$ is the field of moduli of $\mathcal{P}_{w}$ for every $w \in \mathcal{H}$.

This is a simplified form of [S66a, Theorems 5.1 and 6.2], which are applicable to a PEL-type of a more general type. We call $(V, \varphi)$ a canonical model for $\Omega$. We have $K_{\Psi}=\mathbf{Q}$ in Cases SP and UT, and so (1) is trivial in those cases; we will show in Theorem 9.3 below that $k_{\Omega}=k_{N}$ if $\Omega$ is $\Omega^{N}$ of (7.4b). In Case UB we can show that $K_{\Psi}=\mathbf{Q}$ if $m_{v}=n_{v}$ for every $v \in \mathbf{a}$ and $K_{\Psi}$ is a CM-field otherwise; here $m_{v}, n_{v}$ are as in (4.10).

It may be added that no complete proof of the existence of $(V, \varphi)$ with property (4) was given, even when $\Gamma=S p(n, \mathbf{Z})$, in any paper published before 1966.
9.2. Lemma. Let $\sigma$ be an isomorphism of $k_{\Omega}$ onto a subfield $k^{\prime}$ of $\overline{\mathbf{Q}}$ over $K_{\Psi}$. Then there exists a PEL-type $\Omega^{\prime}$ with the following properties:
(1) $\Omega^{\prime}=\left\{K, \Psi, L^{\prime}, \lambda \mathcal{T},\left\{u_{i}^{\prime}\right\}_{i=1}^{s}\right\}$ with the same $K, \Psi, \mathcal{T}$ as $\Omega$, a totally positive element $\lambda$ of $F$, and some $L^{\prime},\left\{u_{i}^{\prime}\right\}_{i=1}^{s}$; we can take $\lambda=1$ in Cases SP and UT.
(2) $k^{\prime}=k_{\Omega^{\prime}}$.
(3) If $\alpha \in \operatorname{Aut}(\mathbf{C}), \alpha=\sigma$ on $k_{\Omega}$, and $\mathcal{P}$ is of type $\Omega$, then $\mathcal{P}^{\alpha}$ is of type $\Omega^{\prime}$.
(4) Let $\left(V^{\prime}, \varphi^{\prime}\right)$ be a canonical model for $\Omega^{\prime}$ and let $\mathcal{F}\left(\Omega^{\prime}\right)=\left\{\mathcal{P}_{z}^{\prime} \mid z \in \mathcal{H}\right\}$. Then there exists a $k^{\prime}$-rational biregular map $f$ of $V^{\sigma}$ onto $V^{\prime}$ with the property that $f\left(\varphi(z)^{\alpha}\right)=\varphi^{\prime}(w)$ whenever $\alpha$ is as in (3) and $\left(\mathcal{P}_{z}\right)^{\alpha}$ is isomorphic to $\mathcal{P}_{w}^{\prime}$.

Proof. Pick any $\mathcal{P}$ of type $\Omega$; then $\mathcal{P}^{\alpha}$ is of type $\Omega^{\prime}$ for some PEL-type $\Omega^{\prime}$. This $\Omega^{\prime}$ does not depend on the choice of $\mathcal{P}$ by virtue of [S66a, Proposition 4.1], as explained in [S66a, p.323, lines $1 \sim 4$ ]. From Theorem 9.1 (2) we see that $k^{\prime}=k_{\Omega^{\prime}}$. For $K \neq F$ assertion (1) was proved in [S64, Proposition 5.2]. Recall that, by virtue of (4.17), $\Psi$ is determined by the signature of $\mathcal{T}_{v}$ for $v \in \mathbf{a}$, and vice versa. If $K=F$, we can always put $\mathcal{T}=\eta_{n}$ as explained in $\S 4.3$, and hence $\mathcal{T}^{\prime}=\eta_{n}$ too.

In the proof of the following theorem we shall determine $\Omega^{\prime}$ in Cases Sp and UT without invoking [S64]. In any case, the space $\mathcal{H}$ is common to $\Omega$ and $\Omega^{\prime}$.

Now, to any structure $\mathcal{P}$ of type $\Omega$ we assign a point $\mathfrak{v}(\mathcal{P})$ of $V$ as follows: by Theorem 4.8 (2), $\mathcal{P}$ is isomorphic to $\mathcal{P}_{z}$ for some $z \in \mathcal{H}$; we then put $\mathfrak{v}(\mathcal{P})=\varphi(z)$. This symbol $\mathfrak{v}$ has the following properties:
(i) $k_{\Omega}(\mathfrak{v}(\mathcal{P}))$ is the field of moduli of $\mathcal{P}$.
(ii) $\mathfrak{v}\left(\mathcal{P}_{1}\right)=\mathfrak{v}\left(\mathcal{P}_{2}\right)$ if and only if $\mathcal{P}_{1}$ is isomorphic to $\mathcal{P}_{2}$.
(iii) If $\mathfrak{p}$ is a $\mathbf{C}$-valued discrete place of a field of rationality for $\mathcal{P}$ of type $\Omega$ such that $\mathfrak{p}(a)=a$ for every $a \in k_{\Omega}$ and $\mathcal{P}_{0}$ is the reduction of $\mathcal{P}$ modulo $\mathfrak{p}$, then $\mathcal{P}_{0}$ is of type $\Omega$ and $\mathfrak{p}(\mathfrak{v}(\mathcal{P}))=\mathfrak{v}\left(\mathcal{P}_{0}\right)$.

Properties (i) and (ii) follow immediately from our definition; (iii) was given in [S66a, Theorem 6.2]. For each $\mathcal{Q}$ of type $\Omega^{\prime}$ we define $\mathfrak{v}^{\prime}(\mathcal{Q}) \in V^{\prime}$ in the same manner. Now take any $\alpha$ as in (3). Then $\mathcal{Q}^{\alpha^{-1}}$ is of type $\Omega$, and so $\mathfrak{v}\left(\mathcal{Q}^{\alpha^{-1}}\right)$ is meaningful. We then define a point $\mathfrak{v}^{\sigma}(\mathcal{Q})$ of $V^{\sigma}$ by $\mathfrak{v}^{\sigma}(\mathcal{Q})=\mathfrak{v}\left(\mathcal{Q}^{\alpha^{-1}}\right)^{\alpha}$. In view of (iii) above, $\mathfrak{v}^{\sigma}(\mathcal{Q})$ is well-defined independently of the choice of $\alpha$. Now by [S66a, Theorem 6.7], ( $V, \mathfrak{v}$ ) can be characterized by properties (ii) and (iii). Since we easily see that $\left(V^{\sigma}, \mathfrak{v}^{\sigma}\right)$ has these properties for $\Omega^{\prime}$, that theorem guaratees a $k^{\prime}$-rational biregular map $f$ of $V^{\sigma}$ onto $V^{\prime}$ such that $f\left(\mathfrak{v}^{\sigma}(\mathcal{Q})\right)=\mathfrak{v}^{\prime}(\mathcal{Q})$ for every $\mathcal{Q}$ of type $\Omega^{\prime}$. If $\left(\mathcal{P}_{z}\right)^{\alpha}$ is isomorphic to $\mathcal{P}_{w}^{\prime}$, then $\mathfrak{v}^{\sigma}\left(\left(\mathcal{P}_{z}\right)^{\alpha}\right)=\mathfrak{v}\left(\mathcal{P}_{z}\right)^{\alpha}=\varphi(z)^{\alpha}$, and so $f\left(\varphi(z)^{\alpha}\right)=\mathfrak{v}^{\prime}\left(\left(\mathcal{P}_{z}\right)^{\alpha}\right)=\mathfrak{v}^{\prime}\left(\mathcal{P}_{w}^{\prime}\right)=\varphi^{\prime}(w)$, which is (4). This completes the proof.
9.3. Theorem. Suppose that $\Omega$ is $\Omega^{N}$ of (7.4b) with $L=\mathfrak{r}_{2 n}^{1}$ in Cases $S P$ and $U T$; let $(V, \varphi)$ and $k_{\Omega}$ be as in Theorem 9.1 and let $k_{N}=\mathbf{Q}(\mathbf{e}(1 / N))$. Then the following assertions hold:
(1) $k_{\Omega}=k_{N}$.
(2) $k_{\Omega}(V) \circ \varphi=\mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$.
(3) There exists a model $(W, \psi)$ of $\Gamma^{N} \backslash \mathcal{H}$ such that $W$ is defined over $\mathbf{Q}$ and $\mathbf{Q}(W) \circ \psi=\mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$.

Proof. Once (1) and (2) are established, Theorem 9.1 (4) implies Theorem 8.12 (for an arbitrary $w$ ). Now, given $\Omega=\Omega^{N}$ as in (7.4b) with $L=\mathfrak{r}_{2 n}^{1}$, for each $s \in \mathbf{Z}_{\mathbf{h}}^{\times}$we put $\Omega_{s}=\left\{K, \Psi, L, \eta_{n},\left\{u_{i} \iota(s)\right\}\right\}$, where $u_{i} \iota(s)$ is an element of $N^{-1} L$ determined modulo $L$ as explained in §8.4; let $\mathcal{F}\left(\Omega_{s}\right)=\left\{\mathcal{P}_{z, s} \mid z \in \mathcal{H}\right\}$. Take a point $z_{0}$ of $\mathcal{H}$ generic for $\mathfrak{K}$ over $\mathbf{Q}_{\mathrm{ab}}$; fixing $s$, define an automorphism $\xi$ of $\mathfrak{K}\left[z_{0}\right]$ by $f\left(z_{0}\right)^{\xi}=f^{\tau(\iota(s))}\left(z_{0}\right)$ for $f \in \mathfrak{K}$. By Lemma 8.6 this means $t\left(a, z_{0}\right)^{\xi}=t\left(a \iota(s), z_{0}\right)$ for every $a \in K_{2 n}^{1}$. By Lemma 7.8, or rather by its proof, $\left(\mathcal{P}_{z_{0}}\right)^{\xi}=\mathcal{P}_{z_{0}, s}$. Let $\left(V_{s}, \varphi_{s}\right)$ be a canonical model for $\Omega_{s}$; put $k_{s}=k_{\Omega_{s}}$ and $\mathfrak{K}_{N}=\mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$. By Theorem 8.12, $\mathfrak{K}_{N}\left[z_{0}\right]$ is the field of moduli of $\mathcal{P}_{z_{0}, s}$, which equals $k_{s}\left(\varphi_{s}\left(z_{0}\right)\right)$. Since $k_{N}$ is algebraically closed in $\mathfrak{K}_{N}$, we see that $k_{s}=k_{N}$, and hence $k_{N}\left(\varphi_{s}\left(z_{0}\right)\right)=\mathfrak{K}_{N}\left[z_{0}\right]$. From this we can conclude that $k_{N}=k_{\Omega}$ and $k_{N}\left(V_{s}\right) \circ \varphi_{s}=\mathfrak{K}_{N}$, which gives (1) $\operatorname{and}(2)$ if we take $s=1$.

Let $\sigma$ be the restriction of $\xi$ to $k_{N}$. Taking $\Omega^{\prime}$ of Lemma 9.2 to be $\Omega_{s}$, we obtain a $k_{N}$-rational biregular map of $V^{\sigma}$ to $V_{s}$. This means that we may assume that $V^{\sigma}=V_{s}$, and

$$
\begin{equation*}
\varphi(z)^{\alpha}=\varphi_{s}(w) \text { if } \alpha \in \operatorname{Aut}(\mathbf{C}), \alpha=\sigma \text { on } k_{N}, \text { and }\left(\mathcal{P}_{z}\right)^{\alpha} \text { is isomorphic to } \mathcal{P}_{w . s} . \tag{9.1}
\end{equation*}
$$ In particular, $\varphi\left(z_{0}\right)^{\xi}=\varphi_{s}\left(z_{0}\right)$.

Clearly $\Omega_{s}$ depends only on $\left\{s_{p}\left(\bmod N \mathbf{Z}_{p}\right)\right\}_{p}$. Writing $\Omega_{\sigma}, \varphi_{\sigma}$, and $\mathcal{P}_{z, \sigma}$ for $\Omega_{s}, \varphi_{s}$, and $\mathcal{P}_{z, s}$, we thus have a canonical model $\left(V^{\sigma}, \varphi_{\sigma}\right)$ for $\Omega_{\sigma}$ for each $\sigma \in$ $\operatorname{Gal}\left(k_{N} / \mathbf{Q}\right)$. Since $\left(V^{\sigma}, \varphi_{\sigma}\right)$ is a model of $\Gamma^{N} \backslash \mathcal{H}$, there exists a biregular map $g_{\sigma}: V \rightarrow V^{\sigma}$ such that $\varphi_{\sigma}=g_{\sigma} \circ \varphi$. Then $k_{N}\left(\varphi\left(z_{0}\right)\right)=\mathfrak{K}_{N}\left[z_{0}\right]=k_{N}\left(\varphi_{\sigma}\left(z_{0}\right)\right)$. Since $\varphi\left(z_{0}\right)$ and $\varphi_{\sigma}\left(z_{0}\right)$ are generic on $V$ and $V^{\sigma}$ over $k_{N}$, there exists a $k_{N}$-rational map of $V$ to $V^{\sigma}$ that sends $\varphi\left(z_{0}\right)$ to $\varphi_{\sigma}\left(z_{0}\right)$. Clearly this map must coincide with $g_{\sigma}$; thus $g_{\sigma}$ is $k_{N}$-rational. Now for each $\sigma \in \operatorname{Gal}\left(k_{N} / \mathbf{Q}\right)$ and $f \in \mathfrak{K}_{N}$ we have a well-defined element $f^{\sigma}$ of $\mathfrak{K}_{N}$, which is the same as $f^{\tau(\iota(s))}$ if $\sigma=[s, \mathbf{Q}]$ on $k_{N}$. Define an automorphism $\alpha$ of $\mathfrak{K}\left[z_{0}\right]$ by $f\left(z_{0}\right)^{\alpha}=f^{\sigma}\left(z_{0}\right)$ for $f \in \mathfrak{K}$. Given $\tau \in \operatorname{Gal}\left(k_{N} / \mathbf{Q}\right)$, define similarly $\beta$ by $f\left(z_{0}\right)^{\beta}=f^{\tau}\left(z_{0}\right)$. Then $\left(\mathcal{P}_{z_{0}}\right)^{\alpha}=\mathcal{P}_{z_{0}, \sigma}$ as observed at the beginning, and similarly $\left(\mathcal{P}_{z_{0}}\right)^{\beta}=\mathcal{P}_{z_{0}, \tau}$ and $\left(\mathcal{P}_{z_{0}}\right)^{\alpha \beta}=\mathcal{P}_{z_{0}, \sigma \tau}$. Therefore $\varphi\left(z_{0}\right)^{\alpha}=\varphi_{\sigma}\left(z_{0}\right), \varphi\left(z_{0}\right)^{\beta}=\varphi_{\tau}\left(z_{0}\right)$, and $\varphi\left(z_{0}\right)^{\alpha \beta}=\varphi_{\sigma \tau}\left(z_{0}\right)$. Thus $\left(\left(g_{\sigma}\right)^{\tau} \circ\right.$ $\left.g_{\tau}\right)\left(\varphi\left(z_{0}\right)\right)=\left(g_{\sigma}\right)^{\tau}\left(\varphi_{\tau}\left(z_{0}\right)\right)=\left(g_{\sigma}\right)^{\tau}\left(\varphi\left(z_{0}\right)^{\beta}\right)=g_{\sigma}\left(\varphi\left(z_{0}\right)\right)^{\beta}=\varphi_{\sigma}\left(z_{0}\right)^{\beta}=\varphi\left(z_{0}\right)^{\alpha \beta}=$ $\varphi_{\sigma \tau}\left(z_{0}\right)=g_{\sigma \tau}\left(\varphi\left(z_{0}\right)\right)$, and so $\left(g_{\sigma}\right)^{\tau} \circ g_{\tau}=g_{\sigma \tau}$, since $\varphi\left(z_{0}\right)$ is generic on $V$ over $k_{N}$. Applying a well-known criterion of [W56, Theorem 3] to $\left\{V^{\sigma}, g_{\sigma}\right\}$, we find a Q-rational variety $W$ and a $k_{N}$-rational biregular map $h$ of $V$ onto $W$ such that $h=h^{\sigma} \circ g_{\sigma}$ for every $\sigma \in \operatorname{Gal}\left(k_{N} / \mathbf{Q}\right)$. Put $\psi=h \circ \varphi$. Clearly $(W, \psi)$ is a model of $\Gamma^{N} \backslash \mathcal{H}$. Now $h\left(\varphi\left(z_{0}\right)\right)^{\alpha}=h^{\sigma}\left(\varphi_{\sigma}\left(z_{0}\right)\right)=\left(h^{\sigma} \circ g_{\sigma}\right)\left(\varphi\left(z_{0}\right)\right)=h\left(\varphi\left(z_{0}\right)\right)$. Thus $h\left(\varphi\left(z_{0}\right)\right)$ is rational over the subfield of $\mathfrak{K}_{N}\left[z_{0}\right]$ fixed by the automorphisms $\alpha$. This subfield corresponds to the subficld of $\mathfrak{K}_{N}=\mathcal{A}_{0}\left(\Gamma^{N}, k_{N}\right)$ fixed by $\operatorname{Gal}\left(k_{N} / \mathbf{Q}\right)$, which is $\mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$. Since $k_{N}\left(h\left(\varphi\left(z_{0}\right)\right)\right)=\mathfrak{K}_{N}\left[z_{0}\right]$, we see that $\mathbf{Q}\left(h\left(\varphi\left(z_{0}\right)\right)\right)$ corresponds to $\mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$. This means that $\mathbf{Q}(W) \circ \psi=\mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$, which completes the proof.
9.4. Given an algebraic number field $M$ of finite degree contained in $\mathbf{C}$, we denote its maximal abelian extension contained in $\mathbf{C}$ by $M_{\mathrm{ab}}$. By class field theory there exists a canonical homomorphism of $M_{\mathbf{A}}^{\times}$onto $\mathrm{Gal}\left(M_{\mathrm{ab}} / M\right)$ whose kernel is the closure of the product of $M^{\times}$and the identity component of $M_{\mathrm{a}}^{\times}$. We denote by $[a, M]$ the element of $\operatorname{Gal}\left(M_{\mathrm{ab}} / M\right)$ which is the image of $a \in M_{\mathbf{A}}^{\times}$. (This includes $[t, \mathbf{Q}]$ of $\S 8.1$ as a special case.)

Let us now briefly recall the notion of a reflex of a CM-type. (For details, see $[\mathrm{S} 98, \S \S 8.3$ and 18.5$]$.) Given a CM-type ( $K, \Phi$ ) in the sense of $\S 2.9$, take a Galois extension $L$ of $\mathbf{Q}$ containing $K$ and take elements $\varphi_{\nu}$ of $G$ so that $\Phi=$ $\left\{\varphi_{1}, \ldots, \varphi_{n}\right\}$. (Thus $n=[K: \mathbf{Q}]$.) Let $K^{*}$ be the field generated over $\mathbf{Q}$ by the elements $\sum_{\nu=1}^{n} x^{\varphi_{\nu}}$ for all $x \in K$. Then it can be shown that $K^{*}$ is a CM-field. Put $G=\operatorname{Gal}(L / \mathbf{Q}), H=\operatorname{Gal}(L / K)$, and $H^{*}=\operatorname{Gal}\left(L / K^{*}\right)$. Then we have

$$
\begin{equation*}
\bigsqcup_{\nu=1}^{n} \varphi_{\nu}^{-1} H=\bigsqcup_{\mu=1}^{m} H^{*} \tau_{\mu} \tag{9.2}
\end{equation*}
$$

with $\tau_{\mu} \in G$, where $m=\left[K^{*}: \mathbf{Q}\right] / 2$. We can show that $\left(K^{*}, \Phi^{*}\right)$ with $\Phi^{*}=\left\{\tau_{\mu}\right\}$ is a CM-type, which we call the reflex of $(K, \Phi)$. The field $K^{*}$ is called the reflex field of $(K, \Phi)$. We can also define a map $g:\left(K^{*}\right)^{\times} \rightarrow K^{\times}$by

$$
\begin{equation*}
g(a)=\prod_{\mu=1}^{m} a^{\tau_{\mu}} \quad\left(a \in\left(K^{*}\right)^{\times}\right) \tag{9.3}
\end{equation*}
$$

The map $g$ can be extended naturally to a homomorphism $\left(K^{*}\right)_{\mathbf{A}}^{\times} \rightarrow K_{\mathbf{A}}^{\times}$. We shall denote this map $g$ also by $\operatorname{det} \Phi^{*}$. Since $\left\{\varphi_{\nu}, \varphi_{\nu} \rho\right\}_{\nu=1}^{n}$ gives the set of all embeddings of $K$ into $\mathbf{C}$, we obtain:

$$
\begin{equation*}
g(a) g(a)^{\rho}=N_{K^{*} / \mathbf{Q}}(a) \quad\left(a \in\left(K^{*}\right)_{\mathbf{A}}^{\times}\right) . \tag{9.4}
\end{equation*}
$$

We now consider a CM-algebra $Y=K_{1} \oplus \cdots \oplus K_{t}$, such that $[Y: K]=2 n$, a $\cdot K$-linear ring-injection $h: Y \rightarrow K_{2 n}^{2 n}$ such that

$$
\begin{equation*}
h\left(a^{\rho}\right) \eta_{n}=\eta_{n} h(a)^{*}, \tag{9.5}
\end{equation*}
$$

and the fixed point $w$ of $h\left(Y^{u}\right)$ as in $\S 4.11$. Define CM-types $\left(K_{i}, \Phi_{i}\right)$ as noted at the end of that subsection. Let ( $K_{i}^{*}, \Phi_{i}^{*}$ ) be the reflex of ( $K_{i}, \Phi_{i}$ ), and $Y^{*}$ the composite field of the $K_{i}^{*}$. (Note: $Y$ is a $K$-algebra, but $Y^{*}$ is a subfield of $\mathbf{C}$.) We then define a map $g:\left(Y^{*}\right)_{\mathbf{A}}^{\times} \rightarrow Y_{\mathbf{A}}^{\times}$by

$$
\begin{equation*}
g(x)=\left(\operatorname{det} \Phi_{i}^{*}\left(N_{Y^{*} / K_{i}^{*}}(x)\right)\right)_{i=1}^{t} . \tag{9.6}
\end{equation*}
$$

9.5. Lemma. The notation being as above, for every $x \in\left(Y^{*}\right)_{\mathbf{A}}^{\times}$we have $\nu(h(g(x)))=g(x) g(x)^{\rho}=N_{Y^{*} / \mathbf{Q}}(x)$ and $\operatorname{det}[h(g(x))]=N_{Y^{*} / \mathbf{Q}}(x)^{n}$. Consequently $h(g(x)) \in\left(G_{0}\right)_{\mathbf{A}+}$.

Proof. That $g(x) g(x)^{\rho}=N_{Y^{*} / \mathbf{Q}}(x)$ follows immediately from (9.4) and (9.6). Also from this and (9.5) we see that $h(g(x)) \in \widetilde{G}_{\mathbf{A}+}$ and $\nu(h(g(x)))=g(x) g(x)^{\rho}$. This proves our lemma in Case SP. To compute det $[h(g(x))]$ in Case UT, take the Galois closure over $\mathbf{Q}$ of the composite of the $K_{i}$ and put $G=\operatorname{Gal}(L / \mathbf{Q}), H_{i}=$ $\operatorname{Gal}\left(L / K_{i}\right), H_{i}^{*}=\operatorname{Gal}\left(L / K_{i}^{*}\right), H^{*}=\operatorname{Gal}\left(L / Y^{*}\right)$, and $J=\operatorname{Gal}(L / K)$. Let $R$ be the group-ring of $G$ over $\mathbf{Z}$; for $x=\sum_{\gamma \in G} c_{\gamma} \gamma \in R$ with $c_{\gamma} \in \mathbf{Z}$ put $x^{\prime}=\sum_{\gamma \in G} c_{\gamma} \gamma^{-1}$; for a subgroup $H$ of $G$ put $[H]=\sum_{\alpha \in H} \alpha$. Take $\varphi_{i \nu}$ and $\tau_{i \mu}$ in $G$ so that $\Phi_{i}=\left\{\varphi_{i \nu}\right\}_{\nu}$ and $\Phi_{i}^{*}=\left\{\tau_{i \mu}\right\}_{\mu}$. Then $\left(\sum_{\nu}\left[H_{i}\right] \varphi_{i \nu}\right)^{\prime}=\sum_{\mu}\left[H_{i}^{*}\right] \tau_{i \mu}$ by (9.2). Since $\Psi$ is equivalent to $n$ times a regular representation of $K$ over $\mathbf{Q}$, we have $n[G]=$ $\sum_{i, \nu}[J] \varphi_{i \nu}$. Given $x \in\left(Y^{*}\right)_{\mathbf{A}}^{\times}$, put $x_{i}=N_{Y^{*} / K_{i}^{*}}(x)$ and $y_{i}=\operatorname{det} \Phi_{i}^{*}\left(x_{i}\right)$. Then $g(x)=\left(y_{i}\right)_{i=1}^{t}$ and $\operatorname{det}[h(g(x))]=\prod_{i=1}^{t} N_{K_{i} / K}\left(y_{i}\right)$, since $h$ is equivalent to a regular representation of $Y$ over $K$. Now $N_{K_{i} / K}\left(y_{i}\right)=\prod_{\gamma \in H_{i} \backslash J} y_{i}^{\gamma}=\prod_{\mu} \prod_{\gamma \in H_{i} \backslash J} x_{i}^{\tau_{i \mu} \gamma}$, and $\sum_{\mu}\left[H_{i}^{*}\right] \tau_{i \mu} \sum_{\gamma \in H_{i} \backslash J} \gamma=\sum_{\nu} \varphi_{i \nu}^{-1}\left[H_{i}\right] \sum_{\gamma \in H_{i} \backslash J} \gamma=\sum_{\nu} \varphi_{i \nu}^{-1}[J]$. Since $H^{*} \subset H_{i}^{*}$, we can write $\sum_{\nu} \varphi_{i \nu}^{-1}[J]=\sum_{\lambda}\left[H^{*}\right] \alpha_{i \lambda}$ with some $\alpha_{i \lambda} \in G$. Then $N_{K_{i} / K}\left(y_{i}\right)=$ $\prod_{\lambda} x^{\alpha_{i \lambda}}$, so that $\operatorname{det}[h(g(x))]=\prod_{i=1}^{t} \prod_{\lambda} x^{\alpha_{i \lambda}}$. Since $\sum_{i, \nu} \varphi_{i \nu}^{-1}[J]=n[G]$, the $\alpha_{i \lambda}$ give $n$ times $H^{*} \backslash G$. Therefeore the last double product is $N_{Y^{*} / \mathbf{Q}}(x)^{n}$, which completes the proof.

Now our main theorem on the reciprocity-law can be stated as follows:
9.6. Theorem. Let $Y, h, w$, and $Y^{*}$ be as in $\S 9.4$ in Case SP or UT. Then for every $f \in \mathfrak{K}$ defined at $w$, the value $f(w)$ belongs to $Y_{\mathrm{ab}}^{*}$. Moreover, if $b \in\left(Y^{*}\right)_{\mathbf{A}}^{\times}$, then $f^{\tau(r)}$ with $r=h\left(g(b)^{-1}\right)$ is finite at $w$ and $f(w)^{\left[b, Y^{*}\right]}=f^{\tau(r)}(w)$. (Notice that $\tau(r)$ is meaningful, since $r \in\left(G_{0}\right)_{\mathbf{A}+}$.)

Proof. Fixing $N$, we use the same symbols as in the proof of Theorem 9.3. Given $b \in\left(Y^{*}\right)_{\mathbf{A}}^{\times}$and $r$ as above, by Lemma 8.3 (1) we can put $r=y \iota(s) \alpha$ with $y \in U^{N}, s \in \mathbf{Z}_{\mathbf{h}}^{\times}$, and $\alpha \in G_{0+}$. By Lemma 9.5, $N_{Y^{*} / \mathbf{Q}}(b)=\nu(r)^{-1}=s \cdot \nu(y \alpha)^{-1}$. Take $\varepsilon \in \operatorname{Aut}\left(\mathbf{C} / Y^{*}\right)$ so that $\varepsilon=\left[b, Y^{*}\right]$ on $Y_{\mathrm{ab}}^{*}$; let $\sigma$ be the restriction of $\varepsilon$ to $k_{N}$. Then $\sigma=[s, \mathbf{Q}]$ on $k_{N}$, and we have a canonical model ( $V^{\sigma}, \varphi_{s}$ ) for $\Omega_{\sigma}=\Omega_{s}$. Now let $f=f_{1} \circ \varphi$ with $f_{1} \in k_{N}(V)$. Then $f \in \mathfrak{K}_{N}$ and $f_{1}^{\sigma}$ is meaningful as an element of $k_{N}\left(V^{\sigma}\right)$. Since $\varphi\left(z_{0}\right)$ is a generic point of $V$ over $k_{N}$, we have $f_{1}^{\sigma}\left(\varphi\left(z_{0}\right)^{\xi}\right)=f_{1}\left(\varphi\left(z_{0}\right)\right)^{\xi}=f\left(z_{0}\right)^{\xi}$. The last quantity is $f^{\tau(\iota(s))}\left(z_{0}\right)$ by the definition of $\xi$. Since $\varphi\left(z_{0}\right)^{\xi}=\varphi_{s}\left(z_{0}\right)$ as noted in (9.1), we obtain $f_{1}^{\sigma}\left(\varphi_{s}\left(z_{0}\right)\right)=f^{\tau(\iota(s))}\left(z_{0}\right)$. Since $z_{0}$ is generic for $\mathfrak{K}$ over $k_{N}$, we have thus

$$
\begin{equation*}
\left(f_{1} \circ \varphi\right)^{\tau(\iota(s))}=f_{1}^{\sigma} \circ \varphi_{s} \quad \text { for every } \quad f_{1} \in k_{N}(V) \tag{9.7}
\end{equation*}
$$

Next we consider the member $\mathcal{P}_{w}=\left(A_{w}, \mathcal{C}_{w}, \iota_{w},\left\{t_{i}\right\}\right)$ of our family as in Theorem 9.3 for the CM-point $w$ in question. Put $\mathcal{Q}=\left(A_{w}, \mathcal{C}_{w}, \iota^{\prime}\right)$ with $\iota^{\prime}$ : $Y \rightarrow \operatorname{End}_{\mathbf{Q}}\left(A_{w}\right)$ as in §4.11. Take $e_{0} \in K_{2 n}^{1}$ so that $K_{2 n}^{1}=e_{0} h(Y)$ and define $q: Y_{\mathbf{a}} \rightarrow\left(\mathbf{C}^{n}\right)^{\mathbf{b}}$ by $q(a)=p_{w}\left(e_{0} h(a)\right)$ for $a \in Y_{\mathbf{a}}$. Then (4.39) shows that $q(c a)={ }^{t} \Phi(c) q(a)$ for $c \in Y$. Take a Z-lattice $\mathfrak{a}$ in $Y$ so that $L=e_{0} h(\mathfrak{a})$. Then $q(\mathfrak{a})=p_{w}(L)$, and $t_{w}$ of (7.2) gives an isomorphism of $\left(\mathbf{C}^{n}\right)^{\mathbf{b}} / q(\mathfrak{a})$ onto $A_{w}$; $t_{i}=t_{w}\left(q\left(v_{i}\right)\right)$ with an element $v_{i} \in Y$ such that $e_{0} h\left(v_{i}\right)=u_{i}$. Take $\zeta \in Y$ so that $\operatorname{Tr}_{K / \mathbf{Q}}\left(e_{0} h(a) \eta_{n} e_{0}^{*}\right)=\operatorname{Tr}_{Y / \mathbf{Q}}(\zeta a)$ for every $a \in Y$. Then from (4.25) and (9.5) we obtain

$$
E_{w}(q(a), q(b))=\operatorname{Tr}_{K / \mathbf{Q}}\left(e_{0} h(a) \eta_{n} h(b)^{*} e_{0}^{*}\right)=\operatorname{Tr}_{K / \mathbf{Q}}\left(e_{0} h\left(a b^{\rho}\right) \eta_{n} e_{0}^{*}\right)=\operatorname{Tr}_{Y / \mathbf{Q}}\left(\zeta a b^{\rho}\right)
$$

Thus $\mathcal{Q}$ is of type $\left\{Y,{ }^{t} \Phi, \mathfrak{a}, \zeta\right\}$ with respect to $t_{w}$ in the sense of [S98, $\left.\S \S 18.4,18.7\right]$. We now apply the main theorem of complex multiplication of abelian varieties, as stated in [ S 98 , Theorem 18.8], to $\mathcal{Q}^{\varepsilon}$. Then $\mathcal{Q}^{\varepsilon}$ is of type $\left\{Y,{ }^{t} \Phi, g(b)^{-1} \mathfrak{a}, \mu \zeta\right\}$ with respect to an isomorphism $\xi^{\prime}$ of $\left(\mathbf{C}^{n}\right)^{\mathbf{b}} / q\left(g(b)^{-1} \mathfrak{a}\right)$ to $\left(A_{w}\right)^{\varepsilon}$, where $\mu$ is the positive rational number such that $\mu \mathbf{Z}=N_{Y^{*} / \mathbf{Q}}(b) \mathbf{Z}$; besides, $t_{w}(q(a))^{\varepsilon}=\xi^{\prime}\left(q\left(g(b)^{-1} a\right)\right)$ for every $a \in Y / \mathfrak{a}$. (The present $q$ is not exactly the same as the map $q$ of [S98, §18.4]. However, the only property of $q$ we need in the proof of [S98, Theorems 18.6 and 18.8] is that $q(c x)=\Phi(c) q(x)$ for $c, x \in Y$. The present ${ }^{t} \Phi$ corresponds to $\Phi$ there. Thus there is no problem with the present q.) Now for $a \in Y / \mathfrak{a}$ we have

$$
q\left(a g(b)^{-1}\right)=p_{w}\left(e_{0} h\left(a g(b)^{-1}\right)\right)=p_{w}\left(e_{0} h(a) r\right)=p_{w}\left(e_{0} h(a) y \iota(s) \alpha\right)
$$

so that $q\left(\mathfrak{a} g(b)^{-1}\right)=p_{w}(L \alpha)$ and $q\left(v_{i} g(b)^{-1}\right)=p_{w}\left(u_{i} \iota(s) \alpha\right)$. Thus $\left(\mathcal{P}_{w}\right)^{\varepsilon}$ is of type $\Omega^{\prime}=\left\{K, \Psi, L \alpha, \mu \eta_{n},\left\{u_{i} \iota(s) \alpha\right\}\right\}$ with respect to $p_{w}$ and $\xi^{\prime}$. In other words, $\left(\mathcal{P}_{w}\right)^{\varepsilon}$ is isomorphic to the member $\mathcal{P}_{w}^{\prime}$ of the family $\mathcal{F}\left(\Omega^{\prime}\right)$ at $w$. Put $w^{\prime}=\alpha(w)$ and $\Lambda={ }^{t} M(\alpha, w)$. By (4.31) we have $p_{w}(x \alpha)=\Lambda p_{w^{\prime}}(x)$ for $x \in\left(K_{\mathrm{a}}\right)_{2 n}^{1}$. Since $N_{Y^{*} / \mathbf{Q}}(b)=s \cdot \nu(y \alpha)^{-1}$, we have $\mu=\nu(\alpha)^{-1}$. Therefore we easily see that $\Lambda$ gives an isomorphism of $\mathcal{P}_{w^{\prime}, s}$ onto $\mathcal{P}_{w}^{\prime}$. Thus $\left(\mathcal{P}_{w}\right)^{\varepsilon}$ is isomorphic to $\mathcal{P}_{w^{\prime}, s}$, so that $\varphi(w)^{\varepsilon}=\varphi_{s}\left(w^{\prime}\right)$ by (9.1). Let $f \in \mathfrak{K}_{N}$; take $f_{1} \in k_{N}(V)$ so that $f=f_{1} \circ \varphi$. Then $f(w)^{\varepsilon}=f_{1}(\varphi(w))^{\varepsilon}=f_{1}^{\sigma}\left(\varphi(w)^{\varepsilon}\right)=f_{1}^{\sigma}\left(\varphi_{s}\left(w^{\prime}\right)\right)=\left(f_{1}^{\sigma} \circ \varphi_{s}\right)\left(w^{\prime}\right)=f^{\tau(\iota(s))}(\alpha w)$ by (9.7). (Here notice that if $f$ is finite at $w$, then $f_{1}$ is finite at $\varphi(w)$, and $f^{\tau(\iota(s))}$ is finite at $\alpha w$.) Since $f^{\tau(y)}=f$ by (8.12), we have $f^{\tau(r)}(w)=f^{\tau(\iota(s) \alpha)}(w)=$ $f^{\tau(\iota(s))}(\alpha w)=f(w)^{\varepsilon}$. Since $\tau(r)$ depends only on $b$, we see that $f(w)^{\varepsilon}$ depends only on the restriction of $\varepsilon$ to $Y_{\mathrm{ab}}^{*}$. Therefore $f(w) \in Y_{\mathrm{ab}}^{*}$, and we obtain the desired equality of our theorem for $f \in \mathfrak{K}_{N}$. Since $N$ is arbitrary, this completes the proof.

In the elliptic modular case we have $G_{0}=\widetilde{G}=G L_{2}(\mathbf{Q})$ and $\mathcal{G}+=\widetilde{G}_{\mathbf{A}+}$, and we take an imaginary quadratic field as $Y$. Then $Y^{*}$ is the isomorphic image of $Y$ in C. Thus the above theorem specialized to that case is exactly the principal result of the classical theory of complex multipication given in [S71, Theorem 6.31].

In this book we consider only canonical models associated with PEL-types introduced in Section 4. Actually we can define canonical models of arithmetic quotients of hermitian symmetric spaces that are not necessarily associated with PEL-types; we can even prove a reciprocity-law similar to that of Theorem 9.6 in such cases. For details the reader is referred to [S67], [S70], and [S98, Section 26].
9.7. Let the notation be as in Theorem 9.1, and let $\kappa \in \mathbf{Z}$. For every nonzero element $h$ of $\mathcal{A}_{\kappa \mathbf{a}}(\Gamma)$ we can speak of its divisor on $V$, which is a divisor on the variety $V$ in the sense of algebraic geometry. We denote it by $\operatorname{div}(h)$. Given a
subfield $k$ of $\mathbf{C}$ containing $k_{\Omega}$ and a $k$-rational divisor $X$ on $V$, we define the standard symbols $\mathcal{L}(X)$ and $\mathcal{L}(X, k)$ of linear systems as follows:

$$
\begin{gather*}
\mathcal{L}(X)=\{f \in \mathbf{C}(V) \mid \operatorname{div}(f) \succ-X\},  \tag{9.8a}\\
\mathcal{L}(X, k)=\mathcal{L}(X) \cap k(V) . \tag{9.8b}
\end{gather*}
$$

It is well-known that $\mathcal{L}(X)=\mathcal{L}(X, k) \otimes_{k} \mathbf{C}$. Also we easily see that $\mathcal{L}(\operatorname{div}(h))$ is $\mathbf{C}$ linearly isomorphic to $\mathcal{M}_{\boldsymbol{\kappa} \mathbf{a}}(\Gamma)$ via the map $f \mapsto f h$. (Here we have to exclude the one-dimensional case that requires the cusp condition. Therefore, strictly speaking, we have to modify the proofs of the following theorems in that case, a task we leave to the reader.)

Hereafter we consider only Cases SP and UT, and denote by ( $V_{N}, \varphi_{N}$ ) a model of $\Gamma^{N} \backslash \mathcal{H}$ such that $\mathbf{Q}\left(V_{N}\right) \circ \varphi_{N}=\mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$ established in Theorem 9.3 above. Then for any subfield $k$ of $\mathbf{C}$ we have $k(V) \circ \varphi_{N}=k \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right) \subset \mathcal{A}_{0}\left(\Gamma^{N}, k\right)$. By Lemma 7.4 we have $k(V) \circ \varphi_{N}=\mathcal{A}_{0}\left(\Gamma^{N}, k\right)$. Therefore we can take $\mathcal{A}_{0}\left(\Gamma^{N}\right)$ and $\mathcal{A}_{0}\left(\Gamma^{N}, k\right)$ in place of $\mathbf{C}(V)$ and $k(V)$ in (9.8a, b).
9.8. Proposition. (1) If $0 \neq h \in \mathcal{A}_{\kappa \mathbf{a}}\left(\Gamma^{N}, D\right)$ with $\kappa \in \mathbf{Z}$ and a subfield $D$ of $\mathbf{C}$, then $\operatorname{div}(h)$ considered on $V_{N}$ is $D$-rational.
(2) There exist a positive integer $\lambda$ and a nonzero element $g \in \mathcal{A}_{\lambda}\left(\Gamma^{1}, \mathbf{Q}\right)$ such that $\operatorname{div}(g)$ considered on $V_{N}$ for every $N$ is $\mathbf{Q}$-rational.

Proof. We prove this only in Case UT; Case SP can be treated in a similar and much simpler way. To prove (2), let $m$ be the complex dimension of $\mathcal{H}$. Take $m$ algebraically independent functions in $\mathcal{A}_{0}\left(\Gamma^{1}, \mathbf{Q}\right)$. Multiplying these by the product of the denominators, we obtain algebraically independent functions $f_{1} / f_{0}, \ldots, f_{m} / f_{0}$ in $\mathcal{A}_{0}\left(\Gamma^{1}, \mathbf{Q}\right)$ with $f_{\nu} \in \mathcal{M}_{\mu \mathbf{a}}\left(\Gamma^{M}, \mathbf{Q}\right)$ with some $M$ and $\mu>0$. Let $z_{a b}^{v}$ be the $(a, b)$-entry of the matrix $z_{v}$ which is the $v$-th component of the variable $z=\left(z_{v}\right)_{v \in \mathbf{a}} \in \mathcal{H}$. Let $h_{p}=f_{p} / f_{0}$ and

$$
r=(2 \pi i)^{-m} \partial\left(h_{1}, \ldots, h_{m}\right) / \partial\left(z_{1}, \ldots, z_{m}\right),
$$

where $z_{1}, \ldots, z_{m}$ are an arbitrarily fixed arrangement of the variables $z_{a b}^{v}$ for all $v \in \mathbf{a}$ and $1 \leq a, b \leq n$. In view of Lemma 3.4 (2), we see that $r \in \mathcal{A}_{2 n \mathbf{a}}\left(\Gamma^{1}\right)$. Now for a function of the form $f(z)=\sum_{h} c(h) \mathbf{e}_{\mathbf{a}}^{n}(h z)$ as in (5.22a) we have

$$
\partial f / \partial z_{a b}^{v}=2 \pi i \sum_{h} c(h) h_{b a}^{v} \mathbf{e}_{\mathbf{a}}^{n}(h z),
$$

where $h_{b a}^{v}$ is the image of $h_{b a}$ under $v$. We have also

$$
\begin{equation*}
(2 \pi i)^{-1} f_{0}^{2} \partial h_{p} / \partial z_{a b}^{v}=(2 \pi i)^{-1}\left(f_{0} \partial f_{p} / \partial z_{a b}^{v}-f_{p} \partial f_{0} / \partial z_{a b}^{v}\right) . \tag{9.9}
\end{equation*}
$$

Therefore we easily see that $f_{0}^{2 m} r \in \mathcal{M}_{l \mathbf{a}}\left(\Gamma^{M}, \Phi\right)$, where $\Phi$ is the Galois closure of $K$ over $\mathbf{Q}$ and $l=2 m \mu+2 n$. Moreover, $\operatorname{div}(r)$, considered on $V_{N}$ for any fixed $N$, is the same as the divisor of $d h_{1} \wedge \cdots \wedge d h_{m}$ on $V_{N}$, which is $\mathbf{Q}$-rational. Now we can view $f_{0}^{2 m} r$ as the determinant of a matrix of size $n^{2} \# \mathbf{a}$ whose ( $v, a, b$ )-th column is $(2 \pi i)^{-1}\left[f_{0}^{2} \partial h_{p} / \partial z_{a b}^{v}\right]_{p=1}^{m}$. Call this column vector $(f ; v, a, b)$. Let $\sigma \in \operatorname{Gal}(\Phi / \mathbf{Q})$. Applying $\sigma$ to (9.9), we observe that

$$
(f ; v, a, b)^{\sigma}= \begin{cases}\left(f ; v^{\prime}, a, b\right) & \text { if } v \sigma=v^{\prime} \text { on } K \\ \left(f ; v^{\prime}, b, a\right) & \text { if } v \sigma=v^{\prime} \rho \text { on } K\end{cases}
$$

where $\rho$ is complex conjugation. Therefore we have $r^{\sigma}= \pm r$, and hence $\operatorname{div}\left(r^{\sigma}\right)=$ $\operatorname{div}\left(d h_{1} \wedge \cdots \wedge d h_{m}\right)=\operatorname{div}(r)$. Then we obtain the desired function $g$ of (2) by
$g=\prod_{\sigma} r^{\sigma}$, where $\sigma$ runs over $\operatorname{Gal}(\Phi / \mathbf{Q})$. Let $0 \neq h \in \mathcal{A}_{\kappa \mathbf{a}}\left(\Gamma^{N}, D\right)$. Take $g$ as above. Then $h^{\lambda} / g^{\kappa} \in \mathcal{A}_{0}\left(\Gamma_{N}, D\right)$ and $\operatorname{div}\left(h^{\lambda} / g^{\kappa}\right)=\lambda \cdot \operatorname{div}(h)-\kappa \cdot \operatorname{div}(g)$. Since this divisor and $\operatorname{div}(g)$ are $D$-rational, $\operatorname{div}(h)$ must be $D$-rational. This proves (1).
9.9. Theorem. Let $0<\kappa \in \mathbf{Z}$ and $0<N \in \mathbf{Z}$; let $D$ be an arbitrary sufield of C. Then the following assertions hold:
(1) $\mathcal{M}_{\kappa \mathbf{a}}\left(\Gamma^{N}\right)=\mathcal{M}_{\kappa \mathbf{a}}\left(\Gamma^{N}, D\right) \otimes_{D} \mathbf{C}$ provided $\mathcal{A}_{\kappa \mathbf{a}}\left(\Gamma^{N}, D\right) \neq\{0\}$.
(2) $\mathcal{M}_{\kappa \mathbf{a}}=\mathcal{M}_{\kappa \mathbf{a}}(\mathbf{Q}) \otimes_{\mathbf{Q}} \mathbf{C}$.
(3) Given $f \in \mathcal{A}_{\kappa \mathbf{a}}$, there exists a finitely generated extension $E$ of $\mathbf{Q}$ such that $f \in \mathcal{A}_{\kappa \mathbf{\kappa}}(E)$.
(4) Let $f \in \mathcal{M}_{\kappa \mathbf{a}}$ (resp. $f \in \mathcal{A}_{\kappa \mathbf{a}}$ ) and $\sigma \in \operatorname{Aut}(\mathbf{C})$. Then $f^{\sigma}$, defined as a formal element in §5.9, is indeed an element of $\mathcal{M}_{\kappa \mathbf{a}}$ (resp. $\mathcal{A}_{\kappa \mathbf{a}}$ ).

Proof. Let $0 \neq h \in \mathcal{A}_{\kappa \mathbf{a}}\left(\Gamma^{N}, D\right)$. By Proposition $9.8(1), \operatorname{div}(h)$ is $D$-rational. Then $f \mapsto h f$ is an isomorphism of $\mathcal{L}(\operatorname{div}(h))$ onto $\mathcal{M}_{\kappa \mathbf{a}}\left(\Gamma^{N}\right)$, and this maps $\mathcal{L}(\operatorname{div}(h), D)$ onto $\mathcal{M}_{\kappa \mathbf{a}}\left(\Gamma^{N}, D\right)$. Since $\mathcal{L}(\operatorname{div}(h))=\mathcal{L}(\operatorname{div}(h), D) \otimes_{D} \mathbf{C}$, we obtain (1). Now by Lemma 6.17 there exists a nonzero element $p$ in $\mathcal{M}_{\kappa \mathbf{a}}(\mathbf{Q})$. Thus $\mathcal{M}_{\kappa \mathbf{a}}\left(\Gamma^{N}, \mathbf{Q}\right) \neq\{0\}$ for sufficiently large $N$. Therefore (2) follows from (1). From (2) we see that every element of $\mathcal{M}_{\kappa \text { a }}$ is a finite $\mathbf{C}$-linear combination of elements of $\mathcal{M}_{\kappa \mathbf{a}}(\mathbf{Q})$. Clearly this implies (4) for $f \in \mathcal{M}_{\kappa \mathbf{a}}$. Given $f \in \mathcal{A}_{\kappa \mathbf{a}}$, we see that $p^{-1} f \in \mathcal{A}_{0}$, so that by Theorem 8.11 (4) and (5), $p^{-1} f \in \mathcal{A}_{0}(E)$ with a finitely generated extension $E$ of $\mathbf{Q}$, and $\left(p^{-1} f\right)^{\sigma} \in \mathcal{A}_{0}$. Then we obtain (3) and (4) for $f \in \mathcal{A}_{\kappa \mathbf{a}}$.
9.10. In (5.1) we defined the symbol $M_{\alpha}(z)=\left(\mu_{v}(\alpha, z)\right)_{v \in \mathbf{b}}$. Let $\omega$ be a representation of $G L_{n}(\mathbf{C})^{\mathbf{b}}$ given in the form $\omega(x)=\bigotimes_{v \in \mathbf{b}} \omega_{v}\left(x_{v}\right)$ with $\mathbf{Q}$-rational representations $\omega_{v}$ of $G L_{n}(\mathbf{C})$. We recall that

$$
\begin{equation*}
\left(f \|_{\omega} \alpha\right)(z)=\omega\left(M_{\alpha}(z)\right)^{-1} f(\alpha z) \tag{9.10}
\end{equation*}
$$

in both Cases SP and UT.
Now, for $\sigma \in \operatorname{Gal}(\overline{\mathbf{Q}} / \mathbf{Q})$, we define a representation $\omega^{\sigma}$ of $G L_{n}(\mathbf{C})^{\mathbf{b}}$ by

$$
\begin{equation*}
\omega^{\sigma}(x)=\bigotimes_{v \in \mathbf{b}} \omega_{v}\left(x_{v \sigma}\right) \tag{9.11}
\end{equation*}
$$

Thus both $\omega$ and $\omega^{\sigma}$ have the same representation space.
9.11. Proposition. For each fixed $v \in \mathbf{b}$ definie a representation $\tau_{v}: G L_{n}(\mathbf{C})^{\mathbf{b}}$ $\rightarrow G L_{n}(\mathbf{C})$ by

$$
\begin{equation*}
\tau_{v}(x)=\operatorname{det}(x)^{\mathbf{b}} x_{v} \quad \text { for } \quad x \in G L_{n}(\mathbf{C})^{\mathbf{b}} \tag{9.12}
\end{equation*}
$$

where $a^{\mathbf{b}}=\prod_{v \in \mathbf{b}} a_{v}$. $\left(\right.$ Thus $\left(\tau_{v}\right)^{\sigma}=\tau_{v \sigma}$.) Given $z_{0} \in \mathcal{H}$, there exist a set of $\mathbf{C}_{n}^{n}$ valued functions $\left\{R_{v}\right\}_{v \in \mathbf{b}}$ and a congruence subgroup $\Gamma$ of $G$ with the following properties:
(1) The columns of $R_{v}$ belong to $\mathcal{M}_{\tau_{v}}\left(\Gamma, K^{v}\right)$, where $K^{v}$ is the image of $K$ under $v$.
(2) $R_{v}^{\sigma}=R_{v \sigma}$ for every $\sigma \in \operatorname{Gal}(\overline{\mathbf{Q}} / \mathbf{Q})$.
(3) $\operatorname{det}\left(R_{v}\left(z_{0}\right)\right) \neq 0$ for every $v \in \mathbf{b}$.
(4) The columns of $R_{v} \|_{\tau_{v}} \alpha$ belong to $\mathcal{M}_{\tau_{v}}\left(K^{v} \mathbf{Q}_{\mathrm{ab}}\right)$ for every $\alpha \in \widetilde{G}_{+}$.

Moreover, these assertions are true with the following modifications: Replace $\mathcal{M}_{\tau_{v}}$ by $\mathcal{A}_{\sigma_{v}}$, where $\sigma_{v}(x)=x_{v}$; (3) should then read: $R_{v}$ is holomorphic at $z_{0}$ and $\operatorname{det}\left(R_{v}\left(z_{0}\right)\right) \neq 0$ for every $v \in \mathbf{b}$.

Proof. We first prove $(1,2,3)$ for $\mathcal{M}_{\tau_{v}}$ in Case SP with $F=\mathbf{Q}$. Put

$$
\begin{equation*}
\psi(u, z ; \lambda)=(2 \pi i)^{-1}\left[\left(\partial / \partial u_{j}\right) \varphi(u, z ; \lambda)\right]_{j=1}^{n} \tag{9.13}
\end{equation*}
$$

where $\varphi$ is $\varphi_{F}$ of (6.17) with $F=\mathbf{Q}$. We view this as a column vector. Take $\Lambda \subset \mathcal{S}\left(\mathbf{Q}_{\mathrm{h}}^{n}\right)$ as in Theorem 6.12 and take $\Gamma_{\lambda}$ as in (6.34). Differentiating (6.34), we obtain

$$
\begin{equation*}
\psi(u, \gamma z ; \lambda)=h_{\gamma}(z) \mu_{\gamma}(z) \psi\left({ }^{t} \mu_{\gamma}(z) u, z ; \lambda\right) \quad \text { for every } \quad \gamma \in \Gamma_{\lambda} . \tag{9.14}
\end{equation*}
$$

Put $\Gamma_{0}=\bigcap_{\lambda \in \Lambda} \Gamma_{\lambda}$. Denoting the elements of $\Lambda$ by $\lambda_{1}, \ldots, \lambda_{m}$, put

$$
\varphi_{k}(u, z)=\varphi\left(u, z ; \lambda_{k}\right), \quad \psi_{k}(u, z)=\psi\left(u, z ; \lambda_{k}\right) \quad(1 \leq k \leq m) .
$$

Since $u \mapsto\left(\varphi_{k}(u, z)\right)_{k=1}^{m}$ is the biregular embedding of Theorem 6.12 (3), we have

$$
\operatorname{rank}\left[\begin{array}{lll}
\varphi_{1}(u, z) & \cdots & \varphi_{m}(u, z) \\
\psi_{1}(u, z) & \cdots & \psi_{m}(u, z)
\end{array}\right]=n+1
$$

for every $(u, z) \in \mathbf{C}^{n} \times \mathfrak{H}$. Therefore, given $z_{0} \in \mathfrak{H}_{n}$, changing the order of $\lambda_{k}$, we may assume that $\operatorname{det}\left[\psi_{1}\left(0, z_{0}\right) \cdots \psi_{n}\left(0, z_{0}\right)\right] \neq 0$. Also we can find an index $j$ such that $\varphi_{j}\left(0, z_{0}\right) \neq 0$. Define a $\mathbf{C}_{n}^{n}$-valued function $R$ by

$$
\begin{equation*}
R(z)=\varphi_{j}(0, z)\left[\psi_{1}(0, z) \cdots \psi_{n}(0, z)\right] . \tag{9.15}
\end{equation*}
$$

Then (6.34) together with (9.14) and Theorem 6.8 (5) shows that $R(\gamma(z))=$ $j_{\gamma}(z) \mu_{\gamma}(z) R(z)$ for every $\gamma$ in a subgroup $\Gamma$ of $\Gamma_{0}$ of finite index. We easily see that $\psi_{k}(0, z)$ has Fourier coefficients in $\mathbf{Q}$. This proves $(1,2,3)$ in the case $F=\mathbf{Q}$.

Next we consider Case SP with $F \neq \mathbf{Q}$. We take the embedding $\varepsilon: \mathfrak{H}_{n}^{\mathbf{a}} \rightarrow \mathfrak{H}_{e n}$ of (6.10), and given $z_{0} \in \mathfrak{H}_{n}^{\text {a }}$, choose the above $R$ on $\mathfrak{H}_{\text {en }}$ (that is, with en instead of $n$ ) so that $\operatorname{det}\left[R\left(\varepsilon\left(z_{0}\right)\right)\right] \neq 0$. With $B$ as in (6.5) put

$$
\left[\begin{array}{c}
R_{1}(z)  \tag{9.16}\\
\vdots \\
R_{e}(z)
\end{array}\right]={ }^{t} B R(\varepsilon(z)) Q \quad\left(z \in \mathfrak{H}_{n}^{\mathrm{a}}\right)
$$

with $Q \in \mathbf{Q}_{n}^{e n}$. Here each $R_{j}$ is an $n \times n$-matrix. We can choose $Q$ so that $\operatorname{det}\left[R_{j}\left(z_{0}\right)\right] \neq 0$ for every $j$. (To see this, we first note an easy fact: Given a nonzero polynomial $p\left(x_{1}, \ldots, x_{m}\right)$, with complex coefficients, there exist rational numbers $q_{1}, \ldots, q_{m}$ such that $p\left(q_{1}, \ldots, q_{m}\right) \neq 0$. Then take a variable $e n \times n$ matrix $X$ and put ${ }^{t} B R\left(\varepsilon\left(z_{0}\right)\right) X=\left[\begin{array}{c}Y_{1} \\ \vdots \\ Y_{e}\end{array}\right]$, where each $Y_{j}$ is square and of size $n$, and apply the above fact to $p(X)=\prod_{j=1}^{e} \operatorname{det}\left(Y_{j}\right)$.) Writing $\left\{R_{v}\right\}_{v \in \mathbf{a}}$ for $\left\{R_{j}\right\}$ and employing ( $6.11 \mathrm{a}, \mathrm{b}$ ), we can easily verify that the $R_{v}$ satisfy $(1,2,3)$ with a suitable choice of $\Gamma \subset S p(n, F)$.

Finally take Case UT. Using the symbols of $\S 6.5$, denote by $T_{v}$ the above function $R_{v}$ with $2 n$ in place of $n$ such that $\operatorname{det}\left[T_{v}\left(\psi\left(z_{0}\right)\right)\right] \neq 0$ for a given $z_{0} \in \mathcal{H}_{n}^{\mathbf{a}}$. Then we define $n \times n$-matrices $S_{v}$ and $R_{v}$ by

$$
\left[\begin{array}{l}
S_{v}(w)  \tag{9.17}\\
R_{v}(w)
\end{array}\right]=A_{v}^{*} T_{v}(\psi(w)) Y \quad\left(w \in \mathcal{H}_{n}^{\mathbf{a}}, v \in \mathbf{a}\right)
$$

with $Y \in \mathbf{Q}_{n}^{2 n}$. By virtue of ( $6.22 \mathrm{~b}, \mathrm{c}$ ) we easily find a congruence subgroup $\Gamma$ of $G$ such that

$$
S_{v}(\alpha w)=j_{\alpha}(w)^{\mathbf{b}} \lambda_{v}(\alpha, w) S_{v}(w) \quad \text { and } \quad R_{v}(\alpha w)=j_{\alpha}(w)^{\mathbf{b}} \mu_{v}(\alpha, w) R_{v}(w)
$$

for every $\alpha \in \Gamma$. Also, with a suitable choice of $Y$, we have $\operatorname{det}\left[S_{v}\left(z_{0}\right) R_{v}\left(z_{0}\right)\right] \neq 0$. Putting $R_{v \rho}=S_{v}$, we can easily verify that the set $\left\{R_{v}\right\}_{v \in \mathbf{b}}$ satisfies (1, 2, 3).

As for (4), since every element of $\widetilde{G}_{+}$is a product of an element of $G$ and an element of the form $\operatorname{diag}\left[1_{n}, c 1_{n}\right]$ with $c \in F, \gg 0$, it is sufficient to treat the case where $\alpha \in G$. Then by means of (6.11a, b), (6.22b, c), (9.16), and (9.17), we can reduce our problem to the case $G=S p(n, \mathbf{Q})$. In this case we prove (4) in a stronger form:
(9.18) If $X \in \mathcal{M}_{\tau}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, then $X \|_{\tau} \alpha \in \mathcal{M}_{\tau}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ for every $\alpha \in S p(n, \mathbf{Q})$.

This is clear if $\alpha \in P \cap G$ with $P$ of Lemma 7.5. Therefore, by that lemma, it is sufficient to prove (9.18) when $\alpha=\eta$. First we take $X$ to be $R$ of (9.15). Let $r(z)=\operatorname{det}(-i z)^{1 / 2}$. By Theorem 6.9 (2) we have

$$
\varphi\left(u, \eta z ; \lambda_{k}\right)=r(z) \varphi\left({ }^{t} \mu_{\eta}(z) u, z ; \lambda_{k}^{\prime}\right)
$$

with $\lambda_{k}^{\prime}$ given by (6.33) with $\lambda_{k}$ as $\lambda$ and $F=\mathbf{Q}$. Clearly $\lambda_{k}^{\prime}$ is $\mathbf{Q}_{\mathrm{ab}}$-valued. Then differentiation shows that $\psi\left(u, \eta z ; \lambda_{k}\right)=r(z) \mu_{\eta}(z) \psi\left({ }^{t} \mu_{\eta}(z) u, z ; \lambda_{k}^{\prime}\right)$. Therefore we obtain

$$
(R \| \eta)(z)=(-i)^{n} \varphi\left(0, z ; \lambda_{j}^{\prime}\right)\left[\psi\left(0, z ; \lambda_{1}^{\prime}\right) \quad \cdots \quad \psi\left(0, z ; \lambda_{n}^{\prime}\right]\right.
$$

This is clearly $\mathbf{Q}_{\mathrm{ab}}$-rational. Now given $X$ as in (9.18), put $Z=R^{-1} X$. Then the entries of $Z$ belong to $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, and so by Theorem 7.10 (4) the entries of $Z \circ \eta$ belong to $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$. Thus $X \| \eta=(R \| \eta)(Z \circ \eta) \in \mathcal{M}_{\tau}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, which proves (9.18). This completes the proof of (4).

As for the assertions for $\mathcal{A}_{\sigma_{v}}$ instead of $\mathcal{M}_{\tau_{v}}$, the desired functions can be obtained by taking $\varphi_{j}(0, z)^{-1}$ in place of $\varphi_{j}(0, z)$ in (9.15). Then the above proof is applicable to this case too, with no other changes.
9.12. Lemma. Given $\gamma=\left[\begin{array}{ll}p & q \\ r & s\end{array}\right] \in S L_{2}(\mathbf{Z})$ with $s>0$, put $\beta=\left[\begin{array}{ll}p 1_{n} & q 1_{n} \\ r 1_{n} & s 1_{n}\end{array}\right]$. Then for some positive integer $M$ we can choose the function $R$ of Proposition 9.11 with $R_{v} \in \mathcal{A}_{\sigma_{v}}$ so that it has the following property: If $\tau \in \operatorname{Gal}(\overline{\mathbf{Q}} / \mathbf{Q}), t \in \mathbf{Z}_{\mathbf{h}}^{\times}, \tau=$ $[t, \mathbf{Q}]$ on $\mathbf{Q}_{\mathrm{ab}}$, and $\gamma-\operatorname{diag}\left[t^{-1}, t\right] \prec M \mathbf{Z}$, then $\left(R_{v} \|_{\sigma_{v}} \eta\right)^{\tau}=R_{v \tau} \|_{\sigma_{v \tau}} \beta \eta$.

Proof. We first consider Case SP with $F=\mathbf{Q}$. Given $\lambda \in \mathcal{S}\left(\mathbf{Q}_{\mathbf{h}}^{n}\right)$ and $t \in \mathbf{Z}_{\mathbf{h}}^{\times}$, define $\lambda_{t}, \lambda^{\prime} \in \mathcal{S}\left(\mathbf{Q}_{\mathbf{h}}^{n}\right)$ by $\lambda_{t}(x)=\lambda(t x)$ for $x \in \mathbf{Q}_{\mathbf{h}}^{n}$ and

$$
\begin{equation*}
\lambda^{\prime}(x)=\int_{\mathbf{Q}_{\mathbf{h}}^{n}} \lambda(y) \mathbf{e}_{\mathbf{h}}\left({ }^{t} x y\right) d y \tag{9.19}
\end{equation*}
$$

as in (6.33). Also, for $\tau \in \operatorname{Gal}(\overline{\mathbf{Q}} / \mathbf{Q})$ define $\lambda^{\tau}$ by $\lambda^{\tau}(x)=\lambda(x)^{\tau}$. Put $\varphi(z, \lambda)=$ $\varphi(0, z ; \lambda)$ and $\psi(z, \lambda)=\psi(0, z ; \lambda)$. Then our construction shows that each column of $R$ is of the form $f(z)=\varphi(z ; \lambda)^{-1} \psi(z, \ell)$ with some $\mathbf{Q}$-valued elements $\lambda, \ell$ of $\mathcal{S}\left(\mathbf{Q}_{\mathrm{h}}^{n}\right)$. By Theorem $6.9(2)$ we have $f \| \eta=\varphi\left(z ; \lambda^{\prime}\right)^{-1} \psi\left(z, \ell^{\prime}\right)$, and hence $(f \| \eta)^{\tau}=$ $\varphi\left(z ;\left(\lambda^{\prime}\right)^{\tau}\right)^{-1} \psi\left(z,\left(\ell^{\prime}\right)^{\tau}\right)$. Take a positive integer $M$ so that

$$
\left\{\xi \in S p(n, \mathbf{Q})_{\mathbf{A}} \mid \xi-1_{2 n} \prec M \mathbf{Z}\right\} \subset D_{\lambda} \cap D_{\ell}
$$

for $D_{\lambda}$ and $D_{\ell}$ of Theorem 6.8 (4). Then, by (6.32), $f \| \beta=\varphi\left(z ; \lambda_{t}\right)^{-1} \psi\left(z, \ell_{t}\right)$. Applying $\tau$ to (9.19), we find that $\left(\lambda^{\prime}\right)^{\tau}=\left(\lambda_{t}\right)^{\prime}$ and $\left(\ell^{\prime}\right)^{\tau}=\left(\ell_{t}\right)^{\prime}$, and hence $(f \| \eta)^{\tau}=\varphi\left(z ;\left(\lambda_{t}\right)^{\prime}\right)^{-1} \psi\left(z,\left(\ell_{t}\right)^{\prime}\right)=f \| \beta \eta$, which proves the case $G=S p(n, \mathbf{Q})$.

To prove the case of $G=S p(n, F)$, we employ the symbols of $\S 6.2$. Let $\beta$ be the element of $G$ defined in $S p(n, F)$ as in our lemma, and let $\beta^{\prime}=\left[\begin{array}{ll}p 1_{e n} & q 1_{e n} \\ r 1_{e n} & s 1_{e n}\end{array}\right] \in$
$S p(e n, \mathbf{Q})$. With $B$ as in (6.5), put $A=B \cdot{ }^{t} B$ and $\alpha=\operatorname{diag}\left[A^{-1}, A\right]$. Let $M^{\prime}$ be the integer with which $R$ on $\operatorname{Sp}(e n, \mathbf{Q})$ has the desired property. We take it so that $R \| \Gamma^{M^{\prime}}=R$. Now take a positive multiple $M$ of $M^{\prime}$ so that $M \alpha \prec M^{\prime} \mathbf{Z}$. With $\omega$ as in (6.7), observe that $\omega\left(\eta_{n}\right)=\eta_{e n} \alpha$ and $\omega(\beta)=\xi \beta^{\prime}$ with

$$
\xi=\left[\begin{array}{cc}
1+q r(1-A) & p q(A-1) \\
s r\left(A^{-1}-1\right) & 1+q r\left(1-A^{-1}\right)
\end{array}\right] .
$$

Therefore $\xi-1 \prec M^{\prime} \mathbf{Z}$ if $q, r \in M \mathbf{Z}$. Combining our result in the case $F=\mathbf{Q}$ with (9.16), we obtain

$$
\begin{equation*}
\left(\left(R_{v} \| \eta\right)^{\tau}\right)_{v \in \mathbf{a}}={ }^{t} B^{\tau}\left(R \| \eta_{e n} \alpha\right)^{\tau}(\varepsilon(z)) Q={ }^{t} B^{\tau}\left(R \| \beta^{\prime} \eta_{e n} \alpha\right)(\varepsilon(z)) Q, \tag{*}
\end{equation*}
$$

since $(g \| \alpha)^{\tau}=g^{\tau} \| \alpha$ holds for $g \in \mathcal{A}_{\sigma}$. On the other hand

$$
\left(R_{v} \| \beta \eta\right)_{v \in \mathbf{a}}={ }^{t} B\left(R \| \xi \beta^{\prime} \eta_{e n} \alpha\right)(\varepsilon(z)) Q={ }^{t} B\left(R \| \beta^{\prime} \eta_{e n} \alpha\right)(\varepsilon(z)) Q
$$

since $R \| \xi=R$. Comparing this with ( ${ }^{*}$ ), we obtain the desired equality for $R_{v}$ on $S p(n, F)$. Case UT can be handled in a similar manner by means of (9.17).
9.13. Theorem. (1) Let $f \in \mathcal{M}_{\omega}$ (resp. $f \in \mathcal{A}_{\omega}$ ) with $\omega$ as in $\S 9.10$ and let $\sigma \in \operatorname{Aut}(\mathbf{C})$. Then $f^{\sigma}$, defined as a formal element in §5.9, is indeed an element of $\mathcal{M}_{\omega^{\sigma}}$ (resp. $\mathcal{A}_{\omega^{\sigma}}$ ). Thus $\left(\mathcal{A}_{\omega}\right)^{\sigma}=\mathcal{A}_{\omega^{\sigma}}$ and $\left(\mathcal{M}_{\omega}\right)^{\sigma}=\mathcal{M}_{\omega^{\sigma}}$.
(2) Given $f \in \mathcal{A}_{\omega}$, there exists a finitely generated extension $k$ of $\mathbf{Q}$ such that $f \in \mathcal{A}_{\omega}(k)$.
(3) $\mathcal{M}_{\omega}(D)$ is stable under $f \mapsto f \|_{\omega} \alpha$ for every $\alpha \in \widetilde{G}_{+}$and every subfield $D$ of $\mathbf{C}$ containing $\mathbf{Q}_{\mathrm{ab}}$ and the Galois closure of $K$ over $\mathbf{Q}$.

Proof. We may assume that $\omega$ is irreducible. Then there is an integer $e$ such that $\omega(c y)=c^{e} \omega(y)$ for $c \in \mathbf{C}^{\times}$. Take $\Gamma$ and $R=\left(R_{v}\right)_{v \in \mathbf{b}}$ as in Proposition 9.11. Then we see that $\omega(R) \circ \alpha=j_{\alpha}^{\text {eb }} \omega\left(\mu_{\alpha} R\right)$ for every $\alpha \in \Gamma$. Take a positive integer $m$ so that $m>e$ and $\operatorname{det}(y)^{m} \omega_{v}(y)^{-1}$ is a polynomial in $y$ for every $v \in \mathbf{b}$; put $s(x)=\operatorname{det}(x)^{-m \mathbf{b}} \omega(x)$ for $x \in G L_{n}(\mathbf{C})^{\mathbf{b}}$. Then $s(R) \circ \alpha=j_{\alpha}^{-\kappa \mathbf{b}} \omega\left(\mu_{\alpha}\right) s(R)$ with $\kappa=m(1+n|\mathbf{b}|)-e$. Given $f \in \mathcal{A}_{\omega}$, put $g=s(R)^{-1} f$, and observe that the components of $g$ belong to $\mathcal{A}_{\kappa \mathbf{b}}$ and that $s(R)^{\sigma}=s\left(R^{\sigma}\right)=s^{\sigma}(R)$. Identify the representation space of $\omega$ with $\mathbf{C}^{t}$ with some $t$ so that $\omega\left(G L_{n}(\mathbf{Q})^{\mathbf{b}}\right)$ acts on $\mathbf{Q}^{t}$. Then $g \in\left(\mathcal{A}_{\kappa \mathbf{b}}\right)^{t}$. By Theorem $9.9(4), g^{\sigma} \in\left(\mathcal{A}_{\kappa \mathbf{b}}\right)^{t}$, so that $f^{\sigma}$, being equal to $s\left(R^{\sigma}\right) g^{\sigma}$, must be defined as an element of $\mathcal{A}_{\omega^{\sigma}}$. Suppose $f \in \mathcal{M}_{\omega}$. Then $g \in\left(\mathcal{M}_{\kappa \mathbf{b}}\right)^{t}$, since $s(R)^{-1}$ is holomorphic everywhere. Now for any point $z_{0}$, take $R$ so that $\operatorname{det}\left[\prod_{v \in \mathbf{b}} R_{v}\left(z_{0}\right)\right] \neq 0$. Now $g^{\sigma} \in\left(\mathcal{M}_{\kappa \mathbf{b}}\right)^{t}$ by Theorem 9.9 (4), so that $f^{\sigma}=s(R)^{\sigma} g^{\sigma}$ is holomorphic at $z_{0}$. Thus $f^{\sigma}$ is holomorphic everywhere, and so $f^{\sigma} \in \mathcal{M}_{\omega^{\sigma}}$. This proves (1). Now $s(R)$ is rational over the Galois closure of $K$. Since $f=s(R) g$, we obtain (2) by applying Theorem 9.9 (3) to $g$. To prove (3), we take $R_{v}$ of Proposition 9.11 with $R_{v} \in \mathcal{A}_{\sigma_{v}}$. Given $f \in \mathcal{M}_{\omega}(D)$, put $g=\omega(R)^{-1} f$. Then $g$ has components in $\mathcal{A}_{0}(D)$. Now $R_{v} \|_{\sigma_{v}} \alpha \in \mathcal{A}_{\sigma_{v}}\left(K^{v} \mathbf{Q}_{\mathrm{ab}}\right)$ as stated in Proposition 9.11, and $g \circ \alpha$ has components in $\mathcal{A}_{0}(D)$ by Theorem 7.10, and hence $f \|_{\omega} \alpha=\left(\omega(R) \|_{\omega} \alpha\right)(g \circ \alpha) \in \mathcal{A}_{\omega}(D)$. This together with (5.30) proves (3).

## 10. Automorphisms of the spaces of automorphic forms

10.1. We are going to prove a theorem analogous to Theorem 8.10 , taking automorphic forms instead of automorphic functions. Let $\omega$ and $\psi$ be two Q-rational
representations of $G L_{n}(\mathbf{C})^{\mathbf{b}}$. For $f \in \mathcal{A}_{\omega}$ and $g \in \mathcal{A}_{\psi}$ we denote by $f \otimes g$ the element of $\mathcal{A}_{\omega \otimes \psi}$ defined in a natural way.

We now define a subgroup $\mathfrak{G}$ of $\mathcal{G}_{+} \times \operatorname{Gal}(\overline{\mathbf{Q}} / \mathbf{Q})$ by

$$
\begin{equation*}
\mathfrak{G}=\left\{(\xi, \sigma) \in \mathcal{G}_{+} \times \operatorname{Gal}(\overline{\mathbf{Q}} / \mathbf{Q}) \mid \sigma(\xi)=\sigma \quad \text { on } \quad \mathbf{Q}_{\mathrm{ab}}\right\} \tag{10.1}
\end{equation*}
$$

where $\sigma(\xi)$ is defined in the paragraph preceding (8.11).
10.2. Theorem. Each element $(\xi, \sigma)$ of $\mathfrak{G}$ gives a $\mathbf{Q}$-linear bijection of $\mathcal{A}_{\omega}(\overline{\mathbf{Q}})$ onto $\mathcal{A}_{\omega^{\sigma}}(\overline{\mathbf{Q}})$, written $f \mapsto f^{(\xi, \sigma)}$, with the following properties:
(1) $\mathcal{M}_{\omega}(\overline{\mathbf{Q}})^{(\xi, \sigma)}=\mathcal{M}_{\omega^{\sigma}}(\overline{\mathbf{Q}})$.
(2) If $f \in \mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, then $f^{(\xi, \sigma)}=f^{\tau(\xi)}$ with $\tau$ of Theorem 8.10.
(3) $(a f)^{(\xi, \sigma)}=a^{\sigma} f^{(\xi, \sigma)}$ for every $a \in \overline{\mathbf{Q}}$.
(4) $\left(f^{(\xi, \sigma)}\right)^{(\zeta, \tau)}=f^{(\xi \zeta, \sigma \tau)}$.
(5) $f^{(\alpha, 1)}=f \|_{\omega} \alpha$ if $\alpha \in \widetilde{G}_{+}$and $f \in \mathcal{A}_{\omega}$.
(6) $(f \otimes g)^{(\xi, \sigma)}=f^{(\xi, \sigma)} \otimes g^{(\xi, \sigma)}$.
(7) $f^{(\xi, \sigma)}$ coincides with $f^{\sigma}$ of $\S 5.9$ and Theorem 9.13 (1) if $\xi=\iota(s)$ with $s \in \mathbf{Z}_{\mathbf{h}}^{\times}$.
(8) If $f \in \mathcal{A}_{\omega}\left(\Gamma^{N}, \overline{\mathbf{Q}}\right)$ and $(\xi, \sigma) \in \mathfrak{G}$ with $\xi \in T^{N}$, then $f^{(\xi, \sigma)}=f^{\sigma}$.

Proof. We first consider the action of $\mathfrak{G}$ on $\mathcal{A}_{0}(\overline{\mathbf{Q}})$. By Lemma 7.4 (1) and Theorem 7.10 (3), $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ and $\overline{\mathbf{Q}}$ are linearly disjoint over $\mathbf{Q}_{\mathrm{ab}}$, and $\mathcal{A}_{0}(\overline{\mathbf{Q}})=$ $\overline{\mathbf{Q}} \mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$. Therefore, given $(\xi, \sigma) \in \mathfrak{G}$, we can find an automorphism of $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ that coincides with $\sigma$ on $\overline{\mathbf{Q}}$ and with $\tau(\xi)$ on $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$. Writing this action by putting $(\xi, \sigma)$ on the upper right, we can easily verify all the assertions restricted to $\mathcal{A}_{0}(\overline{\mathbf{Q}})$. (As for (7) and (8), we can derive them from Theorem 8.10 (3), Theorem 8.11 (2), and (8.13).)

Next, to treat $\mathcal{A}_{\omega}$ with $\omega$ of a general type, let $R=\left(R_{v}\right)_{v \in \mathbf{b}}$ denote the set of functions obtained in Proposition 9.11 with $R_{v} \in \mathcal{A}_{\sigma_{v}}$; take a positive integer $M$ so that the columns of $R_{v}$ and $R_{v} \| \eta$ belong to $\mathcal{A}_{\sigma_{v}}\left(\Gamma^{M}\right)$ for every $v \in \mathbf{b}$ and have the property stated in Lemma 9.12. By Proposition 9.11 (2), the columns of $\omega(R)^{\sigma}$ belong to $\mathcal{A}_{\omega^{\sigma}}\left(\Gamma^{M}\right)$ for every $\sigma \in \operatorname{Gal}(\overline{\mathbf{Q}} / \mathbf{Q})$. Given $f \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$ and $(\xi, \sigma) \in \mathfrak{G}$, by (8.10) we have $\xi \in x \alpha \widetilde{G}_{\mathbf{a}+}$ with $x \in T^{M}$ and $\alpha \in \widetilde{G}_{+}$. Observe that the components of $\omega(R)^{-1} f$ belong to $\mathcal{A}_{0}(\overline{\mathbf{Q}})$, and so $\left(\omega(R)^{-1} f\right)^{(\xi, \sigma)}$ is meaningful. Then we define $f^{(\xi, \sigma)}$ by

$$
\begin{equation*}
f^{(\xi, \sigma)}=\left(\omega^{\sigma}(R) \|_{\omega^{\sigma}} \alpha\right)\left(\omega(R)^{-1} f\right)^{(\xi, \sigma)} \tag{10.2}
\end{equation*}
$$

By Proposition 9.11 (4) the columns of $\omega^{\sigma}(R) \|_{\omega^{\sigma}} \alpha$ belong to $\mathcal{M}_{\omega^{\sigma}}(\overline{\mathbf{Q}})$, and hence $f^{(\xi, \sigma)}$ is indeed an element of $\mathcal{A}_{\omega^{\sigma}}(\overline{\mathbf{Q}})$. Also we can easily verify that this does not depend on the choice of $x$ and $\alpha$. Then clearly (5) and (6) hold. If $\xi=\iota(s)$ with $s \in \mathbf{Z}_{\mathbf{h}}^{\times}$, we can take $x=\iota(s)$ and $\alpha=1$; then we obtain (7). If $\omega$ is the trivial representation, it is consistent with the above action on $\mathcal{A}_{0}$. Also, if $f \in \mathcal{A}_{\omega}\left(\Gamma^{N}, \overline{\mathbf{Q}}\right)$ with a multiple $N$ of $M$ and $x \in T^{N}$, then

$$
\left(\omega(R)^{-1} f\right)^{(\xi, \sigma)}=\left(\omega(R)^{-1} f\right)^{(x, \sigma)(\alpha, 1)}=\left(\omega(R)^{-1} f\right)^{\sigma} \circ \alpha
$$

since (4), (5), and (8) are true for $\mathcal{A}_{0}(\overline{\mathbf{Q}})$, and hence we easily see that

$$
\begin{equation*}
f^{(\xi, \sigma)}=f^{\sigma} \|_{\omega^{\sigma}} \alpha \tag{10.3}
\end{equation*}
$$

This combined with Theorem 9.13 (1) proves (1). We also note that

$$
\begin{equation*}
\omega(R)^{(x \alpha \cdot \sigma)}=\omega^{\sigma}(R) \|_{\omega^{\sigma}} \alpha \quad \text { if } \quad x \in T^{M} \quad \text { and } \quad \alpha \in \widetilde{G}_{+}, \tag{10.4a}
\end{equation*}
$$

which is a special case of (10.2). Also from this and (10.2) with $R_{v}$ as $f$, we obtain

$$
\begin{equation*}
\omega(R)^{(\xi, \sigma)}=\omega\left(\left(R_{v}^{\prime}\right)_{v \in \mathbf{b}}\right) \quad \text { with } \quad R_{v}^{\prime}=R_{v}^{(\xi, \sigma)} \tag{10.4b}
\end{equation*}
$$

Let us now prove (8) assuming (4). Let $f \in \mathcal{A}_{\omega}\left(\Gamma^{N}, \overline{\mathbf{Q}}\right)$ and $(\xi, \sigma) \in \mathfrak{G}$ with $\xi \in T^{N}$. Take a common multiple $N^{\prime}$ of $M$ and $N$; let $\xi \in \beta x \widetilde{G}_{\mathbf{a}+}$ with $\beta \in \widetilde{G}_{+}$and $x \in T^{N^{\prime}}$. Then $\beta \in \Gamma^{N}$ by Lemma 8.3 (2), and so $f^{(\xi, \sigma)}=\left(f^{(\beta, 1)}\right)^{(x, \sigma)}=f^{(x, \sigma)}$. Taking $(x, \sigma)$ and 1 as $(\xi, \sigma)$ and $\alpha$ in (10.3), we obtain $f^{(x, \sigma)}=f^{\sigma}$. This proves (8).

To prove the associativity of (4), we first observe that it follows easily from our definition (10.2) if both $\xi$ and $\zeta$ are contained in $T^{M}$, or $\zeta \in \widetilde{G}_{+}$and $\tau=1$. Now assume that

$$
\begin{equation*}
\left(g^{(\alpha, 1)}\right)^{(\zeta, \tau)}=g^{(\alpha \zeta, \tau)} \quad \text { for every } g \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}}), \alpha \in \widetilde{G}_{+},(\zeta, \tau) \in \mathfrak{G} \tag{}
\end{equation*}
$$

Given $(\xi, \sigma)$ and $(\zeta, \tau)$ in $\mathfrak{G}$, let $\xi \in x \alpha \widetilde{G}_{\mathbf{a}+}$ as above and $\alpha \zeta \in y \beta \widetilde{G}_{\mathbf{a}+}$ with $y \in T^{M}$ and $\beta \in \widetilde{G}_{+}$. Then, assuming ( ${ }^{*}$ ), we have

$$
\begin{aligned}
\left(g^{(\xi, \sigma)}\right)^{(\zeta, \tau)} & =\left(\left(g^{(x, \sigma)}\right)^{(\alpha, 1)}\right)^{(\zeta, \tau)}=\left(g^{(x, \sigma)}\right)^{(\alpha \zeta, \tau)}=\left(g^{(x, \sigma)}\right)^{(y \beta, \tau)} \\
& =\left(\left(g^{(x, \sigma)}\right)^{(y, \tau)}\right)^{(\beta, 1)}=\left(\left(g^{(x y, \sigma \tau)}\right)^{(\beta, 1)}=g^{(x y \beta, \sigma \tau)}=g^{(\xi \zeta, \sigma \tau)},\right.
\end{aligned}
$$

which is the desired equality of (4). Thus our task is to prove (*). Observe that if $\left({ }^{*}\right)$ is true for fixed $(\alpha, g)$ and an arbitrary $\zeta$, then it is true for $\left(\alpha^{-1}, g\right)$ and an arbitrary $\zeta$; if it is true for some fixed $\alpha,(\zeta, \tau)$, and for the columns of $\omega(R)$, then it is true for all $g \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$ and the same $\alpha,(\zeta, \tau)$, since the problem can be reduced to $\omega(R)^{-1} g$, whose components belong to $\mathcal{A}_{0}(\overline{\mathbf{Q}})$. Therefore it is sufficient to prove that

$$
\left(\omega(R)^{(\alpha, 1)}\right)^{(\zeta, \tau)}=\omega(R)^{(\alpha \zeta, \tau)}
$$

for $\alpha$ belonging to a set of generators, say $B$, of $\widetilde{G}_{+}$. Given $\alpha \in B$ and $(\zeta, \tau) \in \mathfrak{G}$, let $\zeta \in y \beta \widetilde{G}_{\mathbf{a}+}$ with $y \in T^{M}$ and $\beta \in \widetilde{G}_{+}$. For simplicity put $S=\omega(R)$. Suppose $\left(S^{(\alpha, 1)}\right)^{(y, \tau)}=S^{(\alpha y, \tau)}$. Then

$$
\left(S^{(\alpha, 1)}\right)^{(\zeta, \tau)}=\left(\left(S^{(\alpha, 1)}\right)^{(y, \tau)}\right)^{(\beta, 1)}=\left(S^{(\alpha y, \tau)}\right)^{(\beta, 1)}=S^{(\alpha y \beta, \tau)}=S^{(\alpha \zeta, \tau)}
$$

Thus it is sufficient to prove

$$
\begin{equation*}
\left(S^{(\alpha, 1)}\right)^{(y, \tau)}=S^{(\alpha y, \tau)} \quad \text { for every } \quad \alpha \in B \quad \text { and } \quad y \in T^{M} \tag{}
\end{equation*}
$$

By Lemma 7.5 we can take $B=\left(P \cap \widetilde{G}_{+}\right) \cup\{\eta\}$. We first consider the case $\alpha=\eta$. Let $y \in T^{M}$. Then $y \in \iota(t) U^{M}$ with $t \in \mathbf{Z}_{\mathbf{h}}^{\times}$. Take $\left[\begin{array}{ll}p & q \\ r & s\end{array}\right] \in S L_{2}(\mathbf{Z})$ so that $s>0$ and $\left[\begin{array}{ll}p & q \\ r & s\end{array}\right]-\left[\begin{array}{cc}t^{-1} & 0 \\ 0 & t\end{array}\right] \prec M \mathbf{Z}$; put $\beta=\left[\begin{array}{cc}p 1_{n} & q 1_{n} \\ r 1_{n} & s 1_{n}\end{array}\right]$ and $w=\eta y \eta^{-1} \beta^{-1}$. Then $w \in T^{M I}$, and $S^{(\eta y \cdot \tau)}=S^{(w \beta \eta \cdot \tau)}=S^{\tau} \|_{\omega^{\tau}} \beta \eta$ by (10.4a). On the other hand, taking $S \| \eta$ as $f$ in (10.2), we have $(S \| \eta)^{(y . \tau)}=S^{\tau} \cdot\left[S^{-1}(S \| \eta)\right]^{(y . \tau)}$. Since $S^{-1}(S \| \eta)$ has components in $\mathcal{A}_{0}\left(\Gamma^{M}\right)$, assertion (8) for $\mathcal{A}_{0}$ shows that $\left[S^{-1}(S \| \eta)\right]^{(y, \tau)}=$ $\left(S^{-1}\right)^{\tau}(S \| \eta)^{\tau}$, so that $(S \| \eta)^{(y \cdot \tau)}=(S \| \eta)^{\tau}$. Thus ( $\left.{ }^{* *}\right)$ with $\alpha=\eta$ can be written

$$
\begin{equation*}
\left(S \|_{\omega} \eta\right)^{\tau}=S^{\tau} \|_{\omega^{\tau}} \beta \eta \tag{***}
\end{equation*}
$$

Since $S=\omega(R)$, the desired equality $\left({ }^{* * *}\right)$ follows immediately from Lemma 9.12.
Next, let $\alpha=\left[\begin{array}{ll}a & b \\ 0 & d\end{array}\right] \in P \cap \widetilde{G}_{+}$. Since $P$ is generated by $P \cap \mathfrak{r}_{2 n}^{2 n}$, we may assume that $\alpha \prec \mathfrak{r}$. Let $N=p M \cdot N_{K / \mathbf{Q}}(\operatorname{det}(\alpha))$ with a positive integer $p$ which will be determined afterward. Given $(y, \tau)$ as in $\left(^{* *}\right)$, let $y \in w \beta \widetilde{G}_{\mathbf{a}+}$ with $w \in T^{N}$ and $\beta \in \widetilde{G}_{+}$. Then $w \in \iota(s) U^{N}$ with $s \in \mathbf{Z}_{\mathbf{h}}^{\times}$. Take a positive integer $r$ so that $s r-1 \prec N \mathbf{Z}$; put $\gamma=\left[\begin{array}{cc}a & r b \\ 0 & d\end{array}\right]$ and $x=\alpha w \gamma^{-1}$. Then we can easily verify that $x \in T^{M}$. Therefore $S^{(\alpha y, \tau)}=S^{(x \gamma \beta, \tau)}=S^{\tau} \| \gamma \beta$ by (10.4a). On the other hand $(S \| \alpha)^{(y, \tau)}=(S \| \alpha)^{(w \beta, \tau)}=(S \| \alpha)^{\tau} \| \beta$ by (10.3), since $S \| \alpha \in \mathcal{A}_{\omega}\left(\Gamma^{N}, \overline{\mathbf{Q}}\right)$. We can put $S=q^{-1} S_{1}$ with $q \in \mathcal{M}_{k}(\overline{\mathbf{Q}}), k \in \mathbf{Z}^{\mathbf{b}}$, and a matrix $S_{1}$ whose columns belong to $\mathcal{M}_{\psi}(\overline{\mathbf{Q}}), \psi(x)=\operatorname{det}(x)^{k} \omega(x)$. Put $q(z)=\sum_{h} a(h) \mathbf{e}_{\mathbf{a}}^{n}(h z)$ with $a(h) \in \overline{\mathbf{Q}}$ and $S_{1}(z)=\sum_{h} c(h) \mathbf{e}_{\mathbf{a}}^{n}(h z)$ with $\overline{\mathbf{Q}}$-rational matrices $c(h)$. We choose $p$ so that $a(h) \neq 0$ or $c(h) \neq 0$ only if $h \prec p^{-1} \mathfrak{r}$. Then $N d^{-1} h b \prec \mathfrak{r}$ for every such $h$, and hence

$$
\left(S_{1} \| \alpha\right)^{\tau}(z)=\omega^{\tau}(d)^{-1} \sum_{h} c(h)^{\tau} \mathbf{e}_{\mathbf{a}}^{n}\left(r d^{-1} h b\right) \mathbf{e}_{\mathbf{a}}^{n}\left(d^{-1} h a z\right)=\left(S_{1}^{\tau} \| \gamma\right)(z)
$$

Similarly $(q \| \alpha)^{\tau}=q^{\tau} \| \gamma$, and hence $(S \| \alpha)^{\tau}=S^{\tau} \| \gamma$. Therefore $(S \| \alpha)^{(y, \tau)}=$ $(S \| \alpha)^{\tau}\left\|\beta=S^{\tau}\right\| \gamma \beta=S^{(\alpha y, \tau)}$, which is $\left({ }^{* *}\right)$ for the present $\alpha$. This completes the proof.
10.3. Lemma. Let $c_{1}, \ldots, c_{m}$ be $m$ elements of $\mathbf{C}$ linearly independent over a subfield $D$ of $\overline{\mathbf{Q}}$. Then there exists a set $\{\sigma\}$ of $m$ automorphisms of $\mathbf{C}$ over $D$ such that $\operatorname{det}\left(c_{\nu}^{\sigma}\right)_{\sigma, \nu} \neq 0$.

Proof. For $\sigma \in \operatorname{Aut}(\mathbf{C} / D)$ let $H_{\sigma}=\left\{x \in \mathbf{C}^{m} \mid \sum_{\nu=1}^{m} c_{\nu}^{\sigma} x_{\nu}=0\right\}$ and let $J$ be the intersection of $H_{\sigma}$ for all such $\sigma$. Since $J$ is a vector subspace of $\mathbf{C}^{m}$ stable under $\operatorname{Aut}(\mathbf{C} / D)$, it is defined over $D$. Then the linear independence of the $c_{\nu}$ shows that $J=\{0\}$. Therefore we can find a set $\{\sigma\}$ of $m$ elements of Aut $(\mathbf{C} / D)$ such that $\bigcap_{\sigma \in\{\sigma\}} H_{\sigma}=\{0\}$. Then $\operatorname{det}\left(c_{\nu}^{\sigma}\right)_{\sigma, \nu} \neq 0$ as desired.
10.4. Theorem. Suppose that $\omega\left(c 1_{n}\right)$ is a scalar matrix for every $c \in F_{\mathbf{a}}^{\times}$ (which is the case if $\omega$ is irreducible). Let $W$ be a subgroup of $\widetilde{G}_{\mathbf{A}+}$ containing an open subgroup of $\left(G_{1}\right)_{\mathbf{A}}$. Suppose that $W \cap \widetilde{G}_{\mathbf{h}}$ is contained in an open compact subgroup of $\widetilde{G}_{\mathbf{h}}$ and $x W x^{-1}=W$ for every $x \in \iota\left(Z_{\mathbf{h}}^{\times}\right)$; let $\Gamma=\widetilde{G} \cap W$. Further, given a subfield $D$ of $\mathbf{C}$ and a character $\chi: \Gamma \rightarrow \mathbf{T}$ of finite order such that $\Gamma^{N} \cap \Gamma \subset \operatorname{Ker}(\chi)$ for some $N$, put

$$
\begin{gather*}
\mathcal{M}_{\omega}(\Gamma, \chi)=\left\{f \in \mathcal{M}_{\omega}|f|_{\omega} \gamma=\chi(\gamma) f \text { for every } \gamma \in \Gamma\right\},  \tag{10.5}\\
\mathcal{M}_{\omega}(\Gamma, D, \chi)=\mathcal{M}_{\omega}(D) \cap \mathcal{M}_{\omega}(\Gamma, \chi) \tag{10.6}
\end{gather*}
$$

Then the following assertions hold:
(1) $\Gamma$ is contained in a congruence subgroup of $\widetilde{G}$, and contains $\Gamma^{N}$ of (7.6) for some $N$.
(2) $\mathcal{M}_{\omega}(\Gamma, \chi)^{\tau}=\mathcal{M}_{\omega^{\tau}}\left(\Gamma, \chi_{\tau}\right)$ for every $\tau \in \operatorname{Aut}(\mathbf{C})$, where $\chi_{\tau}$ is a character of $\Gamma$ of finite order determined by $\chi, \omega$, and $\tau$. If $\Gamma \subset G$ and $\chi$ is trivial, then $\chi_{T}$ is trivial.
(3) $\mathcal{M}_{\omega}(\Gamma, \chi)=\mathcal{M}_{\omega}(\Gamma, \overline{\mathbf{Q}}, \chi) \otimes_{\overline{\mathbf{Q}}} \mathbf{C}$.
(4) Let $\Phi$ be the Galois closure of $K$ over $\mathbf{Q}$ in $\overline{\mathbf{Q}}$. Then $\mathcal{M}_{\omega}(\Gamma)=\mathcal{M}_{\omega}(\Gamma, \Phi) \otimes_{\Phi}$ C provided $\Gamma \subset G$.
(5) Given $k=\left(k_{v}\right)_{v \in \mathbf{b}} \in \mathbf{Z}^{\mathbf{b}}$, put $k^{\sigma}=\left(k_{v}^{\sigma}\right)_{v \in \mathbf{b}}$ with $k_{v}^{\sigma}=k_{v \sigma^{-1}}$ for every $\sigma \in \operatorname{Gal}(\Phi / \mathbf{Q})$; let $\Phi_{k}$ be the subfield of $\Phi$ determined by

$$
\begin{equation*}
\operatorname{Gal}\left(\Phi / \Phi_{k}\right)=\left\{\sigma \in \operatorname{Gal}(\Phi / \mathbf{Q}) \mid k^{\sigma}=k\right\} \tag{10.7}
\end{equation*}
$$

Then $\mathcal{M}_{k}(\Gamma)=\mathcal{M}_{k}\left(\Gamma, \Phi_{k}\right) \otimes_{\Phi_{k}} \mathbf{C}$ provided $\Gamma \subset G$.
Proof. Clearly $\Gamma$ is contained in a congruence subgroup of $\widetilde{G}$. Take $N>2$ so that $U^{N} \cap\left(G_{1}\right)_{\mathbf{A}} \subset W$. Then by Lemma 8.3 (3), $\Gamma^{N} \subset G_{1} \cap W \subset \Gamma$, which proves (1). Changing $N$ suitably, we may assume that $\chi$ is trivial on $\Gamma^{N}$. Let $\tau \in \operatorname{Aut}(\mathbf{C})$ and $f \in \mathcal{M}_{\omega}(\Gamma, \overline{\mathbf{Q}}, \chi)$. Put $x=\iota(s)$ with an element $s$ of $\mathbf{Z}_{\mathbf{h}}^{\times}$such that $[s, \mathbf{Q}]=\tau$ on $\mathbf{Q}_{\mathrm{ab}}$. Given $\gamma \in \Gamma$, we can find $\alpha \in \Gamma$ such that $\alpha \in \gamma G_{1}$ and $x \gamma x^{-1} \in \alpha\left(W \cap U^{N}\right)$. Indeed, since $\gamma^{-1} x \gamma x^{-1} \in\left(G_{1}\right)_{\mathbf{A}}$, by strong approximation we have $\gamma^{-1} x \gamma x^{-1} \in \varepsilon\left(W \cap U^{N}\right)$ with $\varepsilon \in G_{1}$. Then $\gamma \varepsilon$ gives the desired $\alpha$. Now we can define a function $\chi_{s}$ on $\Gamma$ by $\chi_{s}(\gamma)=\chi(\alpha)$ with such an $\alpha$. Indeed, if $\beta$ is another element of $\Gamma$ such that $\beta \in \gamma G_{1}$ and $x \gamma x^{-1} \in \beta\left(W \cap U^{N}\right)$, then $\alpha^{-1} \beta \in U^{N} \cap G_{1}=\Gamma^{N}$, so that $\chi(\alpha)=\chi(\beta)$. Thus $\chi_{s}$ is well-defined. Moreover we can easily show that $\chi_{s}$ is a character of $\Gamma$, since $N$ can be changed for any larger integer. Put $x \gamma=\alpha y$. Then $y \in U^{N} x$. Since $f \in \mathcal{M}_{\omega}\left(\Gamma^{N}, \overline{\mathbf{Q}}\right)$, by Theorem 10.2 (8) we have $f^{(x, \tau)}=f^{(y, \tau)}=f^{\tau}$, so that $f^{\tau} \| \gamma=f^{(x \gamma, \tau)}=f^{(\alpha y, \tau)}=(f \| \alpha)^{(y, \tau)}$. By our assumption on $\omega$, we have $\omega\left(c 1_{n}\right)=c^{m}=\prod_{v \in \mathbf{a}} c_{v}^{m_{v}}$ for every $c \in F_{\mathbf{a}}^{\times}$ with some $m \in \mathbf{Z}^{\mathbf{a}}$. Thus $f \| \alpha=\nu(\alpha)^{-m / 2} f \mid \alpha=\nu(\alpha)^{-m / 2} \chi(\alpha) f$, and hence $f^{\tau} \mid \gamma=\chi_{\tau}(\gamma) f^{\tau}$ with

$$
\begin{equation*}
\chi_{\tau}(\gamma)=\chi_{s}(\gamma)^{\tau} \prod_{v \in \mathbf{a}} \nu(\gamma)_{v \tau}^{m_{v} / 2}\left(\nu(\gamma)_{v}^{-m_{v} / 2}\right)^{\tau} \tag{10.8}
\end{equation*}
$$

Clearly $\chi_{\tau}$ is a character of $\Gamma$ of finite order; if $\Gamma \subset G$ and $\chi$ is trivial, then clearly $\chi_{\tau}$ is trivial. This proves that the left-hand side of the following equality is contained in the right-hand side.

$$
\begin{equation*}
\mathcal{M}_{\omega}(\Gamma, \overline{\mathbf{Q}}, \chi)^{\tau}=\mathcal{M}_{\omega^{\top}}\left(\Gamma, \overline{\mathbf{Q}}, \chi_{\tau}\right) \tag{10.9}
\end{equation*}
$$

Since we easily see that $\left(\chi_{\tau}\right)_{\tau^{-1}}=\chi$, applying $\tau^{-1}$ to the right-hand side, we obtain the opposite inclusion, which proves (10.9).

Taking $\Gamma^{M}$ with an arbitrary $M$ and a trivial character as $\Gamma$ and $\chi$, we see that $\mathcal{M}_{\omega}(\overline{\mathbf{Q}})^{\tau}=\mathcal{M}_{\omega^{\tau}}(\overline{\mathbf{Q}})$. Now we employ the notation of the proof of Theorem 9.13. Given $f \in \mathcal{M}_{\omega}$, put $g=s(R)^{-1} f$ as in that proof. Then $g \in\left(\mathcal{M}_{\kappa \mathbf{b}}\right)^{t}$. By Theorem 9.9 (2) we can put $g=\sum_{\nu=1}^{m} c_{\nu} g_{\nu}$ with $c_{\nu} \in \mathbf{C}$ and $g_{\nu} \in \mathcal{M}_{\kappa \mathbf{b}}(\overline{\mathbf{Q}})^{t}$. Changing $\left\{c_{\nu}\right\}$ for a $\overline{\mathbf{Q}}$-basis of $\sum_{\nu=1}^{m} \overline{\mathbf{Q}} c_{\nu}$, we may assume that the $c_{\nu}$ are linearly independent over $\overline{\mathbf{Q}}$. Put $f_{\nu}=s(R) g_{\nu}$. Then $f_{\nu} \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$ and $f^{\sigma}=\sum_{\nu=1}^{m} c_{\nu}^{\sigma} f_{\nu}$ for every $\sigma \in \operatorname{Aut}(\mathbf{C} / \overline{\mathbf{Q}})$. Now $f^{\sigma} \in \mathcal{M}_{\omega}$, and hence from Lemma 10.3 we see that $f_{\nu} \in \mathcal{M}_{\omega}$. Since $f_{\nu} \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$, we obtain $f_{\nu} \in \mathcal{M}_{\omega}(\overline{\mathbf{Q}})$ by (5.30). This shows that $\mathcal{M}_{\omega}=\mathcal{M}_{\omega}(\overline{\mathbf{Q}}) \otimes_{\overline{\mathbf{Q}}} \mathbf{C}$. Suppose $f \in \mathcal{M}_{\omega}(\Gamma, \chi)$. Then for every $\gamma \in \Gamma$ we have $\chi(\gamma) f=f\left|\gamma=\sum_{\nu=1}^{m} c_{\nu} f_{\nu}\right| \gamma$. Since $f_{\nu} \mid \gamma \in \mathcal{M}_{\omega}(\overline{\mathbf{Q}})$ and the $c_{\nu}$ are linearly independent over $\overline{\mathbf{Q}}$, we obtain $f_{\nu} \mid \gamma=\chi(\gamma) f_{\nu}$ for every $\nu$ and every $\gamma \in \Gamma$. This proves (3). Combining this with (10.9), we obtain (2).

To prove (4), suppose $\Gamma \subset G$; let $f \in \mathcal{M}_{\omega}(\Gamma, \overline{\mathbf{Q}})$. By Theorem 9.13 (2) we can find a finite extension $\Xi$ of $\Phi$ contained in $\overline{\mathbf{Q}}$ such that $f \in \mathcal{A}_{\omega}(\Xi)$. Clearly we may assume that $\Xi$ is a Galois extension of $\Phi$. Then for every $\sigma \in \operatorname{Gal}(\Xi / \Phi)$ we have,
by (2), $f^{\sigma} \in \mathcal{M}_{\omega}(\Gamma, \overline{\mathbf{Q}}) \cap \mathcal{A}_{\omega}(\Xi)=\mathcal{M}_{\omega}(\Gamma, \Xi)$. Therefore $\sum_{\sigma}(a f)^{\sigma} \in \mathcal{M}_{\omega}(\Gamma, \Phi)$ for every $a \in \Xi$, where the sum is taken over all $\sigma \in \operatorname{Gal}(\Xi / \Phi)$. Then we see that $f \in \mathcal{M}_{\omega}(\Gamma, \Phi) \otimes_{\Phi} \Xi$. Combining this with (3), we obtain (4). To prove (5), take $\omega(x)=\operatorname{det}(x)^{k}$ with $k$ as in (5). Since $\omega^{\sigma}=\omega$ if $\sigma=$ id. on $\Phi_{k}$, we have $\mathcal{M}_{\omega}(\Gamma, \Phi)^{\tau}=\mathcal{M}_{\omega}(\Gamma, \Phi)$ for every $\tau \in \operatorname{Aut}\left(\mathbf{C} / \Phi_{k}\right)$. Then $\sum_{\sigma}(a f)^{\sigma} \in \mathcal{M}_{\omega}\left(\Gamma, \Phi_{k}\right)$ for every $f \in \mathcal{M}_{\omega}(\Gamma, \Phi)$ and every $a \in \Phi$, where the sum is taken over all $\sigma \in$ $\operatorname{Gal}\left(\Phi / \Phi_{k}\right)$. This proves that $\mathcal{M}_{\omega}(\Gamma, \Phi)=\mathcal{M}_{\omega}\left(\Gamma, \Phi_{k}\right) \otimes_{\Phi_{k}} \Phi$, which combined with (4) proves (5).

It may be added that $\chi_{s}(\gamma)$ in (10.8) is often $\chi(\gamma)$. Indeed, suppose $\chi(\gamma)$ depends only on $a_{\gamma}$ and $d_{\gamma}$ modulo some $\mathfrak{r}$-ideal. Then taking a suitable $N$ in the above proof, we see that $a_{\alpha}-a_{\gamma} \prec N \mathfrak{r}$ and $d_{\alpha}-d_{\gamma} \prec N \mathfrak{r}$, and therefore $\chi_{s}(\gamma)=\chi(\gamma)$. Thus $\chi_{\tau}(\gamma)=\chi(\gamma)^{\tau}$ for such a $\chi$ if $\nu(\gamma)=1$. For example, take $W=C[\mathfrak{v}, \mathfrak{z}] \cap\left(G_{0}\right)_{\mathbf{A}+}$ with $C[\mathfrak{n}, \mathfrak{z}]$ of (1.17). Here $\mathfrak{y}$ and $\mathfrak{z}$ are $\mathfrak{r}$-ideals such that $\mathfrak{y z} \subset \mathfrak{r}$; we naturally take $m=n$ in (1.17). Taking a character $\varphi$ of $(\mathfrak{r} / \mathfrak{y z})^{\times}$, put $\chi(\gamma)=\varphi\left(\operatorname{det}\left(d_{\gamma}\right)\right)$ for $\gamma \in \Gamma$. Then we have $\chi_{\tau}=\chi^{\tau}$.
10.5. Lemma. Given $\tau \in \operatorname{Aut}(\mathbf{C})$, a positive integer $N$, and $\alpha \in \widetilde{G}_{+}$, there exist two elements $\beta$ and $\gamma$ of $\widetilde{G}_{+}$such that $f^{\tau} \|_{\omega^{\tau}} \alpha=\left(f \|_{\omega} \beta\right)^{\tau}$ and $\left(f \|_{\omega} \alpha\right)^{\tau}=$ $f^{\tau} \|_{\omega^{\tau}} \gamma$ for every $f \in \mathcal{M}_{\omega}\left(\Gamma^{N}\right)$ and every $\mathbf{Q}$-rational representation $\omega$. Moreover, if $\alpha \in G$, then $\beta$ and $\gamma$ can be taken from $G$.

Proof. Decomposing $\omega$ into irreducible representations, we may assume that $\omega$ satisfies the condition of Theorem 10.4. Also, by Theorem 10.4 (4), we may assume that $f \in \mathcal{M}_{\omega}\left(\Gamma^{N}, \overline{\mathbf{Q}}\right)$. Then we may take $\tau \in \operatorname{Aut}(\overline{\mathbf{Q}})$. Take a multiple $M$ of $2 N$ so that $\Gamma^{M} \subset \alpha^{-1} \Gamma^{N} \alpha$. We have $f^{\tau} \in \mathcal{M}_{\omega^{\tau}}\left(\Gamma^{N}\right)$ by Theorem 10.4 (2), and hence $f^{\tau} \| \alpha \in \mathcal{M}_{\omega^{\tau}}\left(\Gamma^{M}\right)$. Let $x=\iota(r)$ with an element $r$ of $\mathbf{Z}_{\mathbf{h}}^{\times}$such that $\tau=[r, \mathbf{Q}]$ on $\mathbf{Q}_{\mathrm{ab}}$. By (8.10) we can put $x \alpha=\beta y \widetilde{G}_{\mathbf{a}+}$ with $\beta \in \widetilde{G}_{+}$and $y \in T^{M}$. Then $f \| \beta=f^{\left(x \alpha y^{-1}, 1\right)}=\left(f^{(x \alpha, \tau)}\right)^{(y, \tau)^{-1}}=\left(f^{\tau} \| \alpha\right)^{\tau^{-1}}$ by Theorem 10.2 (8). Thus $f^{\tau} \| \alpha=(f \| \beta)^{\tau}$ as desired. Similarly we can put $\alpha x \in z \gamma \widetilde{G}_{\mathbf{a}+}$ with $z \in T^{M}$ and $\gamma \in \widetilde{G}_{+}$. By Theorem $10.2(8),(f \| \alpha)^{(x, \tau)}=(f \| \alpha)^{\tau}$ and $f^{(z, \tau)}=f^{\tau}$, so that $(f \| \alpha)^{\tau}=f^{(\alpha x, \tau)}=f^{(z \gamma, \tau)}=f^{\tau} \| \gamma$ as expected. The last assertion is clear from our choice of $\beta$ and $\gamma$.

Returning to questions (Q1), (Q2), and (Q3) of §5.9, Theorems 9.13 and 10.4 answer (Q1) and (Q2); the above lemma answers (Q3).
10.6. To treat forms of half-integral weight, we naturally confine our discussion to Case SP. Thus $G=S p(n, F)$. By a quasi-representation of $G L_{n}(\mathbf{C})^{\mathbf{a}}$ we understand a symbol $\psi$ given by

$$
\begin{equation*}
\psi(x)=\operatorname{det}(x)^{\mathbf{a} / 2} \omega(x) \tag{10.10}
\end{equation*}
$$

with a representation $\omega$ of $G L_{n}(\mathbf{C})^{\mathbf{a}}$ as in $\S 9.10$. Given $\gamma$ in the group $\Gamma^{\theta}$ of (6.30) and a function $f$ on $\mathfrak{H}_{n}^{\text {a }}$ with values in the representation space of $\omega$, we put

$$
\begin{equation*}
\left(f \|_{\psi} \gamma\right)(z)=h_{\gamma}(z)^{-1}\left(f \|_{\omega} \gamma\right)(z) \tag{10.11}
\end{equation*}
$$

with $h_{\gamma}$ of Theorem 6.8. For a congruence subgroup $\Gamma$ of $G$ contained in $\Gamma^{\theta}$ we define $\mathcal{M}_{\psi}(\Gamma)$ by conditions (5.8), (5.9), and (5.10), taking $f \|_{\psi} \gamma$ instead of $\left.f\right|_{\omega} \gamma$ in (5.9); we then denote by $\mathcal{M}_{\psi}$ the union of $\mathcal{M}_{\psi}(\Gamma)$ for all such $\Gamma$ 's. If $f \in \mathcal{M}_{\psi}$, we have (5.20) and (5.21) with suitable $M$ and $U$, in view of Theorem 6.8 (3).

Therefore we have an expansion of type (5.22a, b) for $f$. Also, Proposition 5.7 is valid for the present $f$, since what is needed in the proof is (5.20) and (5.21). Now for $\sigma \in \operatorname{Aut}(\mathbf{C})$ we can define $f^{\sigma}$ as a formal series by (5.29). In the following theorem we shall prove that $f^{\sigma}$ defines an element of $\mathcal{M}_{\psi^{\sigma}}$ with a certain quasirepresentation $\psi^{\sigma}$. Given $\alpha \in \widetilde{G}_{+}$, let $p(z)$ be any branch of the square root of $j_{\alpha}(z)^{\mathbf{a}}$. Then by Theorem 6.9 (1) we can show that $p(z)^{-1} f \|_{\omega} \alpha \in \mathcal{M}_{\psi}$, and so we have an expansion

$$
\begin{equation*}
p(z)^{-1}\left(f \|_{\omega} \alpha\right)(z)=\sum_{h \in S} c_{\alpha, p}(h) \mathbf{e}_{\mathbf{a}}^{n}(h z) . \tag{10.12}
\end{equation*}
$$

We call $f$ a cusp form if $c_{\alpha, p}(h)=0$ for every $(\alpha, p)$ and for every $h$ such that $\operatorname{det}(h)=0$, and denote by $\mathcal{S}_{\psi}(\Gamma)$ (resp. $\mathcal{S}_{\psi}$ ) the set of all cusp forms contained in $\mathcal{M}_{\psi}(\Gamma)\left(\right.$ resp. $\left.\mathcal{M}_{\psi}\right)$. We can restrict $\alpha$ to $S p(n, F)$ by virtue of Lemma 1.3 (3).

Further, given a subfield $D$ of $\mathbf{C}$ and a character $\chi: \Gamma \rightarrow \mathbf{T}$ of finite order, we define $\mathcal{M}_{\psi}(D), \mathcal{S}_{\psi}(D), \mathcal{A}_{\psi}(D), \mathcal{M}_{\psi}(\Gamma, D), \mathcal{S}_{\psi}(\Gamma, D)$, and $\mathcal{A}_{\psi}(\Gamma, D)$ in the same manner as in §5.8, and put

$$
\begin{gathered}
\mathcal{M}_{\psi}(\Gamma, \chi)=\left\{f \in \mathcal{M}_{\psi} \mid f \|_{\psi} \gamma=\chi(\gamma) f \text { for every } \gamma \in \Gamma\right\} \\
\\
\mathcal{M}_{\psi}(\Gamma, D, \chi)=\mathcal{M}_{\psi}(D) \cap \mathcal{M}_{\psi}(\Gamma, \chi)
\end{gathered}
$$

We put also $\mathcal{A}_{\psi}(\Gamma)=\mathcal{A}_{\psi}(\Gamma, \mathbf{C})$ and $\mathcal{A}_{\psi}=\mathcal{A}_{\psi}(\mathbf{C})$.
Let $k$ be a half-integral weight and let $m=\left(m_{v}\right)_{v \in \mathbf{a}}$ with $m_{v}=k_{v}-1 / 2$ as in §6.10. If $\omega(x)=\operatorname{det}(x)^{m}$, then $f \|_{\psi} \gamma$ coincides with $f \|_{k} \gamma$ of (6.36), and hence the symbol $\mathcal{M}_{\psi}$ and $\mathcal{A}_{\psi}$ coincide with $\mathcal{M}_{k}$ and $\mathcal{A}_{k}$ of $\S 6.10$. We write $\mathcal{S}_{k}$ for $\mathcal{S}_{\psi}$ in this case.
10.7. Theorem (Case SP). Let $\Gamma$ and $\chi$ be as in Theorem 10.4, and $\psi$ be as in (10.10); let $\Phi$ be the Galois closure of $F$ over $\mathbf{Q}$ in $\overline{\mathbf{Q}}$. Suppose that $\Gamma \subset \Gamma^{\theta}$ and that $b_{\gamma} \prec 2 \mathfrak{d}^{-1}$ and $c_{\gamma} \prec 2 \mathfrak{d}$ for every $\gamma \in \Gamma$. Then the following assertions hold:
(1) Given $f \in \mathcal{A}_{\psi}$, there exists a finitely generated extension $D$ of $\mathbf{Q}$ such that $f \in \mathcal{A}_{\psi}(D)$.
(2) $\left(\mathcal{M}_{\psi}\right)^{\tau}=\mathcal{M}_{\psi^{\tau}},\left(\mathcal{A}_{\psi}\right)^{\tau}=\mathcal{A}_{\psi^{\tau}}$, and $\mathcal{M}_{\psi}(\Gamma, \chi)^{\tau}=\mathcal{M}_{\psi^{\tau}}\left(\Gamma, \chi_{\tau}\right)$ for every $\tau \in \operatorname{Aut}(\mathbf{C})$, where $\psi^{\tau}$ is defined by $\psi^{\tau}(x)=\operatorname{det}(x)^{\mathbf{a} / 2} \omega^{\tau}(x)$, and $\chi_{\tau}$ is a character of $\Gamma$ of finite order determined by $\chi, \omega$, and $\tau$ as in Theorem 10.4 (2).
(3) $\mathcal{M}_{\psi}(\Gamma, \chi)=\mathcal{M}_{\psi}(\Gamma, \overline{\mathbf{Q}}, \chi) \otimes_{\overline{\mathbf{Q}}} \mathbf{C}$.
(4) $\mathcal{M}_{\psi}(\Gamma)=\mathcal{M}_{\psi}(\Gamma, \Phi) \otimes_{\Phi} \mathbf{C}$.
(5) Given a half-integral weight $k=\left(k_{v}\right)_{v \in \mathbf{a}} \in 2^{-1} \mathbf{Z}^{\mathbf{a}}$, put $k^{\sigma}=\left(k_{v}^{\sigma}\right)_{v \in \mathbf{a}}$ with $k_{v}^{\sigma}=k_{v \sigma^{-1}}$ for every $\sigma \in \operatorname{Gal}(\Phi / \mathbf{Q})$; let $\Phi_{k}$ be the subfield of $\Phi$ such that

$$
\operatorname{Gal}\left(\Phi / \Phi_{k}\right)=\left\{\sigma \in \operatorname{Gal}(\Phi / \mathbf{Q}) \mid k^{\sigma}=k\right\}
$$

Then $\mathcal{M}_{k}(\Gamma)=\mathcal{M}_{k}\left(\Gamma, \Phi_{k}\right) \otimes_{\Phi_{k}} \mathbf{C}$.
(6) Let $D$ be a subfield of $\mathbf{C}$ containing $\mathbf{Q}_{\mathrm{ab}}$ and $\Phi$. Given $\alpha \in \widetilde{G}_{+}$, let $p(z)$ be any branch of the square root of $j_{\alpha}(z)^{\mathbf{a}}$. For $f \in \mathcal{M}_{\psi}(D)$ put $g(z)=p(z)^{-1} \cdot\left(f \|_{\omega} \alpha\right)$. Then $g \in \mathcal{M}_{\psi}(D)$.
(7) Let $\alpha \in \widetilde{G}_{+}$and $\tau \in \operatorname{Aut}(\mathbf{C})$; let $p(z)$ be as in (6). Then there exists an element $\beta$ of $\widetilde{G}_{+}$and a branch $q(z)$ of the square root of $j_{\beta}(z)^{\text {a }}$ such that $\left(q(z)^{-1}\left(f \|_{\omega} \beta\right)\right)^{\tau}=p(z)^{-1} \cdot\left(f^{\tau} \|_{\omega^{\tau}} \alpha\right)$ for every $f \in \mathcal{M}_{\psi}(\Gamma)$. Moreover, if $\alpha \in G$, then $\beta$ can be taken from $G$.

Proof. Put $\theta(z)=\sum_{a \in \mathfrak{g}^{n}} \mathbf{e}_{\mathbf{a}}\left({ }^{t} a z a / 2\right)$ for $z \in \mathfrak{H}_{n}^{\mathbf{a}}$ and $\zeta(x)=\operatorname{det}(x)^{\mathbf{a}} \omega(x)$. Define a character $\varphi$ of $\Gamma^{\theta}$ by $h_{\gamma}(z)^{2}=\varphi(\gamma) j_{\gamma}^{\text {a }}$ for $\gamma \in \Gamma^{\theta}$. By Theorem 6.8 (5), $\varphi(\gamma)=\prod_{v \mid 2} \varepsilon_{v}\left(\operatorname{det}\left(d_{\gamma}\right)\right)$ if $\gamma \in \Gamma$ with $\varepsilon$ defined there. Let $\tau \in \operatorname{Aut}(\mathbf{C})$. If $f \in \mathcal{M}_{\psi}(\Gamma, \chi)$, then we easily see that $\theta f \in \mathcal{M}_{\zeta}(\Gamma, \varphi \chi)$, so that $(\theta f)^{\tau} \in$ $\mathcal{M}_{\zeta^{\tau}}\left(\Gamma,(\varphi \chi)_{\tau}\right)$ by Theorem 10.4 (2); observe that $(\varphi \chi)_{\tau}=\varphi \chi_{\tau}$. (Notice that we get the same $\chi_{\tau}$ for both $\zeta$ and $\omega$, since $\Gamma \subset G$.) We define a vector-valued meromorphic function $f^{\prime}$ on $\mathfrak{H}_{n}^{\text {a }}$ by $f^{\prime}=\theta^{-1}(\theta f)^{\tau}$. Then $f^{\prime} \|_{\psi^{\tau}} \gamma=\chi_{\tau}(\gamma) f^{\prime}$ for every $\gamma \in \Gamma$. Now $f \otimes f \in \mathcal{M}_{\rho}$ with $\rho(x)=\operatorname{det}(x)^{\mathbf{a}}(\omega \otimes \omega)(x)$, and $f^{\prime} \otimes f^{\prime}=$ $\theta^{-2}(\theta f \otimes \theta f)^{\tau}=(f \otimes f)^{\tau}$. Therefore $f^{\prime} \otimes f^{\prime}$ is holomorphic everywhere, and hence the square of any component of $f^{\prime}$ is holomorphic everywhere. Thus $f^{\prime}$ has a Fourier expansion, which, multiplied by $\theta$, equals $(\theta f)^{\tau}=\theta f^{\tau}$. Since $\theta$ is not a zero-divisor in the ring of formal series of $\S 5.9$, we see that $f^{\tau}$ gives the Fourier expansion of $f^{\prime}$. This proves that $f^{\tau} \in \mathcal{M}_{\psi^{\tau}}\left(\Gamma, \chi_{\tau}\right)$. Considering the action of $\tau^{-1}$ in the same way, we obtain the last equality of (2), which clearly implies the first two equalities. Next, by Theorem 10.4 (3), $\theta f=\sum_{\nu=1}^{m} c_{\nu} g_{\nu}$ with $c_{\nu} \in \mathbf{C}$ and $g_{\nu} \in \mathcal{M}_{\zeta}(\Gamma, \overline{\mathbf{Q}}, \varphi \chi)$. Changing $\left\{c_{\nu}\right\}$ for a $\overline{\mathbf{Q}}$-basis of $\sum_{\nu=1}^{m} \overline{\mathbf{Q}} c_{\nu}$, we may assume that the $c_{\nu}$ are linearly independent over $\overline{\mathbf{Q}}$. Now $f^{\sigma}=\sum_{\nu=1}^{m} c_{\nu}^{\sigma} \theta^{-1} g_{\nu}$ for every $\sigma \in \operatorname{Aut}(\mathbf{C} / \overline{\mathbf{Q}})$. Since $f^{\sigma}$ is holomorphic everywhere, from Lemma 10.3 we see that $\theta^{-1} g_{\nu}$ is holomorphic everywhere. Thus $\theta^{-1} g_{\nu} \in \mathcal{M}_{\psi}(\Gamma, \overline{\mathbf{Q}}, \chi)$. This proves (3). Assertions (4) and (5) can be proved by the same technique as in the proof of Theorem 10.4 (4) and (5).

Next, let the notation be as in (6). Since $\theta f \in \mathcal{M}_{\zeta}(D)$, we have $(\theta f) \|_{\zeta} \alpha \in$ $\mathcal{M}_{\zeta}(D)$ by Theorem 9.13 (3). Now $(\theta f) \|_{\zeta} \alpha=p^{-1}(\theta \circ \alpha) g$. By Theorem 7.11, $p^{-1}(\theta \circ \alpha) \in \mathcal{M}_{\mathbf{a} / 2}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, and hence $g \in \mathcal{A}_{\psi}(D)$. Clearly $g$ is holomorphic, so that $g \in \mathcal{M}_{\psi}(D)$. This proves (6). To prove (7), take $N>2$ so that $\theta^{2} \in \mathcal{M}_{\mathbf{a}}\left(\Gamma^{N}\right), \Gamma^{N} \subset$ $\Gamma$, and $\theta f \in \mathcal{M}_{\zeta}\left(\Gamma^{N}\right)$ for every $f \in \mathcal{M}_{\psi}(\Gamma)$. Fix such an $f$; given $\alpha \in \widetilde{G}_{+}$, take $\beta$ as in Lemma 10.5; take any branch $q(z)$ of the square root of $j_{\beta}(z)^{\text {a }}$. We have $\theta^{2} \| \alpha=\left(\theta^{2} \| \beta\right)^{\tau}$ and $(\theta f)^{\tau} \|_{\varsigma^{\tau}} \alpha=\left((\theta f) \|_{\zeta} \beta\right)^{\tau}$, and so we see that $p^{-1}\left(f^{\tau} \|_{\omega^{\top}} \alpha\right)=$ $\pm\left(q^{-1}\left(f \|_{\omega} \beta\right)\right)^{\top}$. Now let $f$ and $g$ be two elements of $\mathcal{M}_{\psi}(\Gamma)$ linearly independent over $\mathbf{Q}$. Then $p^{-1}\left(\left(a f^{\tau}+b g^{\tau}\right) \|_{\omega^{\tau}} \alpha\right)=\varepsilon_{a, b}\left(q^{-1}\left((a f+b g) \|_{\omega} \beta\right)\right)^{\tau}$ with $\varepsilon_{a, b}= \pm 1$ for every $a, b \in \mathbf{Q}$. By an elementary argument we easily see that $\varepsilon_{a, b}$ is a constant. Changing $q$ accordingly, we obtain (7). Finally, to prove (1), let $f=g^{-1} h$ with $g \in \mathcal{M}_{e}, e \in \mathbf{Z}^{\mathbf{a}}$ and $h \in \mathcal{M}_{\tau_{e}}$, where $\tau_{e}(x)=\operatorname{det}(x)^{e} \psi(x)$. Taking $\tau_{e}$ to be $\psi$ in (4), we see that (1) is true for $h$; it is also true for $g$ by Theorem 9.13 (2). Therefore we obtain (1) for $f$.
10.8. Theorem. Let $\omega$ be as in $\S 9.10$ in Cases SP and UT; let $\psi$ be as in (10.10) in Case SP; let $\Phi$ be the Galois closure of $K$ over $\mathbf{Q}$ in $\mathbf{C}$ in both cases. Then the following assertions hold:
(1) $\left(\mathcal{S}_{\omega}\right)^{\tau}=\mathcal{S}_{\omega^{\top}}$ and $\left(\mathcal{S}_{\psi}\right)^{\tau}=\mathcal{S}_{\psi^{\tau}}$ for every $\tau \in \operatorname{Aut}(\mathbf{C})$.
(2) $\mathcal{S}_{\omega}(\Gamma)=\mathcal{S}_{\omega}(\Gamma, \Phi) \otimes_{\Phi} \mathbf{C}$ and $\mathcal{S}_{\psi}(\Gamma)=\mathcal{S}_{\psi}(\Gamma, \Phi) \otimes_{\Phi} \mathbf{C}$, where $\Gamma$ is a group as in Theorem 10.4 (4) or Theorem 10.7 (4).

Proof. Let $f \in \mathcal{S}_{\omega}$. Given $\alpha \in \widetilde{G}_{+}$and $\tau \in \operatorname{Aut}(\mathbf{C})$, take $\beta$ as in Lemma 10.5. Applying $\tau$ to the Fourier expansion of $f \| \beta$, we find that $f^{\tau}$ is a cusp form. From this we obtain $\left(\mathcal{S}_{\omega}\right)^{\tau}=\mathcal{S}_{\omega^{\top}}$. The same type of reasoning applies to $\mathcal{S}_{\psi}$, if we take $\beta$ as in Theorem 10.7 (7). Thus we obtain (1). To prove (2), take $f \in$ $\mathcal{S}_{\psi}(\Gamma)$. By Theorem 10.7 (4), $f=\sum_{\nu=1}^{m} c_{\nu} g_{\nu}$ with $c_{\nu} \in \mathbf{C}$ and $g_{\nu} \in \mathcal{M}_{\psi}(\Gamma, \overline{\mathbf{Q}})$.

We may assume that the $c_{\nu}$ are linearly independent over $\mathbf{C}$. Take a set of $m$ automorphisms $\{\sigma\}$ of $\mathbf{C}$ over $\overline{\mathbf{Q}}$ as in Lemma 10.3. Then $\sum_{\nu=1}^{m} c_{\nu}^{\sigma} g_{\nu}=f^{\sigma} \in \mathcal{S}_{\psi}(\Gamma)$ by (1), and so $g_{\nu} \in \mathcal{S}_{\psi}(\Gamma)$. This proves that $\mathcal{S}_{\psi}(\Gamma)=\mathcal{S}_{\psi}(\Gamma, \overline{\mathbf{Q}}) \otimes_{\overline{\mathbf{Q}}} \mathbf{C}$. Then we can prove that $\mathcal{S}_{\psi}(\Gamma, \overline{\mathbf{Q}})=\mathcal{S}_{\psi}(\Gamma, \Phi) \otimes_{\Phi} \overline{\mathbf{Q}}$ by a Galois-theoretical argument as in the proof of Theorem 10.4 (4). The same technique applies to $\mathcal{S}_{\omega}$. Thus we obtain (2).
10.9. Theorem. The symbols $Y, Y^{*}, w, b$, and $r$ being the same as in Theorem 9.6, let $\sigma$ be an element of $\operatorname{Gal}(\overline{\mathbf{Q}} / \mathbf{Q})$ such that $\sigma=\left[b, Y^{*}\right]$ on $Y_{\mathrm{ab}}^{*}$. Then the following assertions hold:
(1) If an element $f$ of $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ is finite at $w$, then $f^{(r, \sigma)}$ is finite at $w$ and $f^{(r, \sigma)}(w)=f(w)^{\sigma}$. (This is a generalization of Theorem 9.6.)
(2) If an element $f$ of $\mathcal{A}_{\omega}(\overline{\mathbf{Q}})$ is finite at $w$, then $f^{(r, \sigma)}$ is finite at $w$.
(3) Let $Q$ be a square matrix whose columns belong to $\mathcal{A}_{\omega}(\overline{\mathbf{Q}})$. If $Q$ is finite and invertible at $w$, then $Q^{(r, \sigma)}$ is finite and invertible at $w$.

Proof. To prove (1), let $f \in \mathcal{A}_{0}\left(\Gamma^{N}, \overline{\mathbf{Q}}\right)$, and put $f=g \circ \varphi_{N}$ with $g \in \overline{\mathbf{Q}}\left(V_{N}\right)$ with $\left(V_{N}, \varphi_{N}\right)$ of $\S 9.7$. Let $z_{0}$ be a point of $\mathcal{H}$ generic for $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ over $\overline{\mathbf{Q}}$. Then $\varphi_{N}\left(z_{0}\right)$ is a generic point of $V_{N}$ over $\overline{\mathbf{Q}}$, and so if $f$ is finite at $w$, then $g$ is defined at $\varphi_{N}(w)$, so that $g=p\left(\varphi_{N}\left(z_{0}\right)\right) / q\left(\varphi_{N}\left(z_{0}\right)\right)$ with polynomials $p$ and $q$ such that $q\left(\varphi_{N}(w)\right) \neq 0$. Writing $p$ and $q$ as $\overline{\mathbf{Q}}$-linear combinations of $\mathbf{Q}$-rational polynomials, we find that $f=\sum_{i} a_{i} t_{i} /\left(\sum_{j} b_{j} s_{j}\right)$ with $a_{i}, b_{j} \in \overline{\mathbf{Q}}$ and $t_{i}, s_{j} \in \mathcal{A}_{0}(\mathbf{Q})$ finite at $w$ such that $\sum_{j} b_{j} s_{j}(w) \neq 0$. Then $f^{(r, \sigma)}=\sum_{i} a_{i}^{\sigma} i_{i}^{\tau(r)} /\left(\sum_{j} b_{j}^{\sigma} s_{j}^{\tau(r)}\right)$. By Theorem 9.6, $t_{i}^{\tau(r)}$ and $s_{j}^{\tau(r)}$ are finite at $w$; moreover $\sum_{i} a_{i}^{\sigma} t_{i}^{\tau(r)}(w)=\left(\sum_{i} a_{i} t_{i}(w)\right)^{\sigma}$ and $\sum_{j} b_{j}^{\sigma} s_{j}^{\tau(r)}(w)=\left(\sum_{j} b_{j} s_{j}(w)\right)^{\sigma} \neq 0$. Therefore we obtain (1).

Let us now prove
(10.13) If $f \in \mathcal{A}_{\kappa \mathbf{a}}(\overline{\mathbf{Q}})$ with $\kappa \in \mathbf{Z}$ and $f$ is finite at $w$, then $f^{(r, \sigma)}$ is finite at $w$; moreover, if $f(w) \neq 0$, then $f^{(r, \sigma)}(w) \neq 0$.
This follows from (1) if $\kappa=0$. Put $g(z)=\theta_{K}(0, r z ; \lambda)$ with $0<r \in \mathbf{Q}$ and a Q-valued $\lambda$. By Proposition 6.14, $g \in \mathcal{M}_{\mathbf{a} / 2}(\mathbf{Q})$ in Case SP and $g \in \mathcal{M}_{\mathbf{a}}(\mathbf{Q})$ in Case UT. Choosing suitable $r$ and $\lambda$, we may assume that $g(w) \neq 0$, by virtue of Theorem 6.12 (2). Assuming $\kappa>0$, put $h=g^{2 \kappa}$ in Case SP and $h=g^{\kappa}$ in Case UT. If $f$ belongs to $\mathcal{A}_{\kappa \mathbf{a}}(\overline{\mathbf{Q}})$ and is finite at $w$, then $f / h$ belongs to $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ and is finite at $w$, and hence $f^{(r, \sigma)} / h^{(r, \sigma)}$ is finite at $w$ by (1); therefore $f^{(r, \sigma)}$ is finite at $w$. Suppose $f(w) \neq 0$; put $p=h^{(r, \sigma)^{-1}} / f$. Then $p$ belongs to $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ and is finite at $w$. We have $p^{(r, \sigma)}(w) f^{(r, \sigma)}(w)=h(w) \neq 0$. By (1), $p^{(r, \sigma)}$ is finite at $w$, and hence $f^{(r, \sigma)}(w) \neq 0$. This proves (10.13) for $\kappa>0$. Next let $f \in \mathcal{A}_{-\kappa \mathbf{a}}(\overline{\mathbf{Q}}), \kappa>0$; suppose $f$ is finite at $w$. If $f(w) \neq 0$, then applying our result to $f^{-1}$, we obtain the desired fact. If $f(w)=0$, then apply the last result to $f+h^{-1}$. In this way we obtain (10.13) for every $\kappa \in \mathbf{Z}$.

To prove (2) and (3), take $R_{v} \in \mathcal{M}_{\sigma_{v}}(\overline{\mathbf{Q}})$ as in Proposition 9.11 so that $R_{v}$ is finite at $w$ and $\operatorname{det}\left(R_{v}(w)\right) \neq 0$ for every $v \in \mathbf{b}$; put $q=\prod_{v \in \mathbf{b}} \operatorname{det}\left(R_{v}\right)$. Then $q \in \mathcal{A}_{\mathbf{b}}(\overline{\mathbf{Q}})$. Given $f \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$ finite at $w$, put $g=\omega(R)^{-1} f$. Then $g$ has components in $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ and is finite at $w$. We have $f^{(r, \sigma)}=\omega(R)^{(r, \sigma)} g^{(r, \sigma)}$, and $g^{(r, \sigma)}$ is finite at $w$ by (1). Since $q(w) \neq 0$, we have $q^{(r, \sigma)}(w) \neq 0$ by (10.13). Therefore $\operatorname{det}\left(R_{v}\right)^{(r . \sigma)}(w) \neq 0$, so that by (10.4b), $\omega(R)^{(r . \sigma)}$ is finite at $w$. Thus $f^{(r, \sigma)}$ is finite at $w$, which proves (2).

Let $Q$ be as in (3); suppose $Q$ is finite and invertible at $w$. Observe that the columns of ${ }^{t} Q^{-1}$ belongs to $\mathcal{A}_{\rho}(\overline{\mathbf{Q}})$ with $\rho={ }^{t} \omega^{-1}$. By (2) both $Q^{(r, \sigma)}$ and $\left({ }^{t} Q^{-1}\right)^{(r, \sigma)}$ are finite at $w$. Therefore $Q^{(r, \sigma)}$ is invertible at $w$. This completes the proof.

Some results, similar to, but somewhat different from, Theorem 10.9, for $G=$ $S P(n, \mathbf{Q})$ were given in [S78b, Section 2, Theorem 3.12, and Proposition 3.13].

We insert here an easy lemma which is necessary for our later purposes.
10.10. Lemma. Let $\Phi$ be as in Theorem 10.4 (4), and $K^{\prime}$ be the reflex field defined for $(K, \tau)$ of §3.5 as in §1.12. Let $\omega$ and $\psi$ be as in §9.10 and (10.10). Given $\alpha \in \tilde{G}_{+}$, take $p(z)$ as in Theorem 10.7 (6) and put $\left(f \|_{\psi}(\alpha, p)\right)(z)=$ $p(z)^{-1}\left(f \|_{\omega} \alpha\right)(z)$; let $\sigma \in \operatorname{Aut}\left(\mathbf{C} / \mathbf{Q}_{\mathrm{ab}} \Phi\right)$. Then

$$
\begin{aligned}
\left(f \|_{\omega} \alpha\right)^{\sigma} & =f^{\sigma} \|_{\omega} \alpha \text { for every } f \in \mathcal{M}_{\omega} \\
\left(f \|_{\psi}(\alpha, p)\right)^{\sigma} & =f^{\sigma} \|_{\psi}(\alpha, p) \text { for every } f \in \mathcal{M}_{\psi}
\end{aligned}
$$

Moreover, if $\omega(x)=\operatorname{det}(x)^{\kappa \mathbf{a}}$ with $\kappa \in \mathbf{Z}$, then these hold for every $\sigma \in \operatorname{Aut}\left(\mathbf{C} / \mathbf{Q}_{\mathrm{ab}}\right)$ in Case SP and $\sigma \in \operatorname{Aut}\left(\mathbf{C} / \mathbf{Q}_{\mathrm{ab}} K^{\prime}\right)$ in Case UT.

Proof. Let $D=\mathbf{Q}_{\mathrm{ab}}$ in Case SP and $D=\mathbf{Q}_{\mathrm{ab}} K^{\prime}$ in Case UT if $\omega(x)=$ $\operatorname{det}(x)^{\kappa \mathrm{a}}$; let $D=\mathbf{Q}_{\mathrm{ab}} \Phi$ in the general case. The desired facts for $f$ in $\mathcal{M}_{\omega}(D)$ or $\mathcal{M}_{\psi}(D)$ follow from Theorem 7.11, Theorem 9.13 (3), and Theorem 10.7 (6). Then, for an arbitrary $f$ in $\mathcal{M}_{\omega}$ or $\mathcal{M}_{\psi}$, we obtain our assertion from Theorem 10.4 (5) and Theorem 10.7 (5).

## 11. Arithmeticity at CM-points

11.1. Let us now consider, in all three cases $S P$, $U T$, and $U B$, the family $\mathcal{F}(\Omega)=\left\{\mathcal{P}_{z} \mid z \in \mathcal{H}\right\}, \mathcal{P}_{z}=\left(A_{z}, \mathcal{C}_{z}, \iota_{z} ;\left\{t_{i}(z)\right\}_{i=1}^{s}\right)$, of (4.26) with a PEL-type $\Omega=\left\{K, \Psi, L, \mathcal{T},\left\{u_{i}\right\}_{i=1}^{s}\right\}$ of (4.7). With $\Gamma$ as in (4.28), let $(V, \varphi)$ be a model of $\Gamma \backslash \mathcal{H}$ as in Theorem 9.1. As explained in $\S 5.4$, we can identify $\mathcal{A}_{0}(\Gamma)$ with the function field of $V$. Let $\mathcal{A}_{0}(\Gamma, \overline{\mathbf{Q}})$ denote the set of all functions of the form $f \circ \varphi$ with a $\overline{\mathbf{Q}}$-rational function $f$ on $V$, which is meaningful since $V$ is defined over $\overline{\mathbf{Q}}$. Then let $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ denote the union of the fields $\mathcal{A}_{0}(\Gamma, \overline{\mathbf{Q}})$ for all such $\Gamma$. Clearly $\mathcal{A}_{0}(\Gamma)=\mathbf{C} \mathcal{A}_{0}(\Gamma, \overline{\mathbf{Q}})$ and $\mathcal{A}_{0}=\mathbf{C} \mathcal{A}_{0}(\overline{\mathbf{Q}})$. For an arbitrary congruence subgroup $\Gamma^{\prime}$ of $G$, we put $\mathcal{A}_{0}\left(\Gamma^{\prime}, \overline{\mathbf{Q}}\right)=\mathcal{A}_{0}(\overline{\mathbf{Q}}) \cap \mathcal{A}_{0}\left(\Gamma^{\prime}\right)$. For each $w \in \mathcal{H}$, let $\mathfrak{F}_{w}$ denote the field generated over $\mathbf{Q}$ by $f(w)$ for all $f$ in $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ finite at $w$. From Theorem 9.1 (4) and Theorem 2.8 (2) we see that $\mathfrak{F}_{w}$ contains the field of moduli of $\mathcal{P}_{w}$, which is algebraic over the field of moduli of $\left(A_{w}, \mathcal{C}_{w}\right)$. By Theorem 2.8 (3) we can find a model of $\mathcal{P}_{w}$ rational over $\mathfrak{F}_{w}$. In this section, whenever we speak of $\mathcal{P}_{w}$, we take it to be rational over $\mathfrak{F}_{w}$.

In Cases SP and UT we already defined $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ in $\S 5.8$ in terms of the Fourier coefficients of automorphic forms. That this coincides with $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ defined above can be seen from Theorem 9.3, since $\mathcal{A}_{0}(\overline{\mathbf{Q}})=\bigcup_{N=1}^{\infty} \mathcal{A}_{0}\left(\Gamma^{N}, \overline{\mathbf{Q}}\right)$ and $\mathcal{A}_{0}\left(\Gamma^{N}, \overline{\mathbf{Q}}\right)=$ $\overline{\mathbf{Q}} \mathcal{A}_{0}\left(\Gamma^{N}, \mathbf{Q}\right)$ by Theorem 8.11 (2). We can also speak of the $\overline{\mathbf{Q}}$-rationality of automorphic forms in Cases SP and UT. The main purpose of this section is to define such in Case UB, or rather in all three cases, in terms of a certain property at each CM-point.

Let us hereafter denote by $\mathcal{H}_{\mathrm{CM}}$ the set of all CM-points of $\mathcal{H}$. The field of moduli of $\left(A_{w}, \mathcal{C}_{w}\right)$ for $w \in \mathcal{H}_{\mathrm{CM}}$ is contained in $\overline{\mathbf{Q}}$ as shown in [S98, Proposition

26 on p. 96 or Corollary 18.9]. Therefore if $w \in \mathcal{H}_{\mathrm{CM}}$, then $k_{\Omega}(\varphi(w)) \subset \overline{\mathbf{Q}}$; hence if an element $f$ of $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ is finite at $w$, then $f(w) \in \overline{\mathbf{Q}}$, so that $\mathfrak{F}_{w} \subset \overline{\mathbf{Q}}$. (In Cases SP and UT, Theorem 9.6 gives a more precise result, but we do not need it in this section.) Thus $\mathcal{P}_{w}$ is $\overline{\mathbf{Q}}$-rational for every $w \in \mathcal{H}_{\mathrm{CM}}$.
11.2. Proposition. (1) If $f \in \mathcal{A}_{0}(\overline{\mathbf{Q}})$ and $\alpha \in \widetilde{G}_{+}$, then $f \circ \alpha \in \mathcal{A}_{0}(\overline{\mathbf{Q}})$.
(2) Let $\mathcal{W}$ be a subset of $\mathcal{H}_{\mathrm{CM}}$ dense in $\mathcal{H}$. If $f \in \mathcal{A}_{0}$ and $f(w) \in \overline{\mathbf{Q}}$ for every $w \in \mathcal{W}$ where $f$ is finite, then $f \in \mathcal{A}_{0}(\overline{\mathbf{Q}})$.

Proof. Given $f \in \mathcal{A}_{0}$, we can find $\Omega$ so that $f \in \mathcal{A}_{0}(\Gamma)$ for $\Gamma$ of (4.28). Then $f=g \circ \varphi$ with $g \in \mathbf{C}(V)$. Let $\mathcal{W}^{\prime}$ be the subset of $\mathcal{W}$ consisting of the points where $f$ is defined. Take a field of rationality $k$ for $g$ containing $\overline{\mathbf{Q}}$ and take also an isomorphism $\sigma$ of $k$ onto a subfield of $\mathbf{C}$ over $\overline{\mathbf{Q}}$. Suppose $f(w) \in \overline{\mathbf{Q}}$ for every $w \in \mathcal{W}^{\prime}$. Now $\varphi(w)$ is $\overline{\mathbf{Q}}$-rational, and hence $g^{\sigma}(\varphi(w))=g(\varphi(w))^{\sigma}=$ $f(w)^{\sigma}=f(w)=g(\varphi(w))$. Since $\varphi\left(\mathcal{W}^{\prime}\right)$ is dense in $V$, we obtain $g^{\sigma}=g$, that is, $g$ is $\overline{\mathbf{Q}}$-rational, and so $f \in \mathcal{A}_{0}(\overline{\mathbf{Q}})$. This proves (2). Clearly (1) follows from (2), since an element of $\widetilde{G}_{+}$maps $\mathcal{H}_{\mathrm{CM}}$ into $\mathcal{H}_{\mathrm{CM}}$ as noted in $\S 4.11$.

It should of course be remarked that $\mathcal{W}$ as in (2) exists. Indeed, we have at least one CM-point $w_{0}$ in $\mathcal{H}$, as remarked in $\S 4.11$. Now $G_{\mathbf{a}}$ acts on $\mathcal{H}$ transitively, and the projection of $G$ to $G_{\mathbf{a}}$ is dense in $G_{\mathbf{a}}$. (The last fact in Cases SP and UT is proved in Lemma 7.5. In Case UB, it can be proved by means of the Cayley parametrization; see [S98, Proposition 23.5].) Then we can take the set of points $\alpha\left(w_{0}\right)$ for all $\alpha \in G$ as $\mathcal{W}$.
11.3. Given a CM-field $K$, we denote by $J_{K}$ the set of all embeddings of $K$ into $\mathbf{C}$, and by $\mathrm{id}_{K}$ the identity embedding of $K$ into $\mathbf{C}$. We then denote by $I_{K}$ the free $\mathbf{Z}$-module generated by the elements of $J_{K}$. If $(K, \Phi)$ is a CM-type, then we naturally view $\Phi$ as a subset of $J_{K}$, and denote by the same letter $\Phi$ the element of $I_{K}$ that is the sum of the members of $\Phi$. To make our exposition easier, we assume that every CM-field in this section is a subfield of $\mathbf{C}$. We always denote complex conjugation by $\rho$.

Let us now recall some basic properties of the period symbol associated with a CM-field $K$, which is a bilinear map $p_{K}: J_{K} \times J_{K} \rightarrow \mathbf{C}^{\times} / \overline{\mathbf{Q}}^{\times}$with certain properties (see [S98, Theorem 32.5]). Here we note a fundamental property:

Given a CM-type ( $K, \Phi$ ), take $(A, \iota$ ) of type $(K, \Phi)$ in the sense of $\S 2.9$, rational over $\overline{\mathbf{Q}}$. Let $n=\operatorname{dim}(A)$ and $\Phi=\sum_{\nu=1}^{n} \varphi_{\nu}$ with $\varphi_{\nu} \in J_{K}$. Then for each $\nu$ there exists a holomorphic 1-form $\omega$ on $A$ such that $\delta \iota(a) \omega=a^{\varphi_{\nu}} \omega$ if $a \in K$ and $\iota(a) \in \operatorname{End}(A)$, where $\delta \iota(a)$ denotes the action of $\iota(a)$ on the space of 1-forms. Clearly $\mathbf{C} \omega$ depends only on $\varphi_{\nu}$. In this setting we have (see [S98, Theorem 32.2])
(11.1) $\omega$ is $\overline{\mathbf{Q}}$-rational if and only if the value of the integral $\int_{c} \omega$ belongs to the coset $\pi p_{K}\left(\varphi_{\nu}, \Phi\right) \overline{\mathbf{Q}}$ for every 1-cycle $c$ on $A$ with coefficients in $\mathbf{Z}$.

Take a CM-algebra $Y=K_{1} \oplus \cdots \oplus K_{t}$ with CM-fields $K_{i}$. We denote by $J_{Y}$ the set of all nontrivial ring-homomorphisms of $Y$ into $\mathbf{C}$, and identify $J_{Y}$ with the union $\bigcup_{i=1}^{t} J_{K_{i}}$ in an obvious way. We also denote by $I_{Y}$ the free $\mathbf{Z}$-module generated by the elements of $J_{Y}$, and identify $I_{Y^{\prime}}$ with the direct sum $I_{K_{1}} \oplus \cdots \oplus I_{K_{t}}$. Given $\alpha=\left(\alpha_{i}\right)_{i=1}^{t}$ and $\beta=\left(\beta_{i}\right)_{i=1}^{t}$ in $I_{Y}$ with $\alpha_{i}, \beta_{i} \in I_{K_{i}}$, we define an element
$p_{Y}(\alpha, \beta)$ of $\mathbf{C}^{\times} / \overline{\mathbf{Q}}^{\times}$by

$$
\begin{equation*}
p_{Y}(\alpha, \beta)=\prod_{i=1}^{t} p_{K_{i}}\left(\alpha_{i}, \beta_{i}\right) \tag{11.2}
\end{equation*}
$$

In the following treatment we denote any nonzero complex number belonging to the coset $p_{Y}(\alpha, \beta)$ by the same symbol $p_{Y}(\alpha, \beta)$. The same convention applies also to the symbols $\mathfrak{p}_{v}(w), \mathfrak{P}_{\omega}(w)$, and $\mathfrak{P}_{k}(w)$ which will be defined below.
11.4. Returning to the family $\mathcal{F}(\Omega)$ of $\S 11.1$, we take a CM-algebra $Y$, a map $h: Y \rightarrow K_{r}^{r}$, the fixed point $w$ of $h\left(Y^{u}\right)$, and the injection $\iota^{\prime}: Y \rightarrow \operatorname{End}_{\mathbf{Q}}\left(A_{w}\right)$ as in $\S 4.11$. We have seen there that $\left(A_{w}, \iota^{\prime}\right)$ is of type $(Y, \Phi)$ with some $\Phi$ whose restriction to $K$ is equivalent to $\Psi$. We take here a $\overline{\mathbf{Q}}$-rational model of $\mathcal{P}_{w}$ as explained in §11.1. Let $\Phi_{v}, \psi_{v}$, and $\varphi_{v}$ be as in (4.38) and (4.40) for the present $Y$. For $\alpha \in Y^{u}$ we have $\psi_{v}(\alpha)=\lambda_{v}(h(\alpha), w)$ and $\varphi_{v}(\alpha)=\mu_{v}(h(\alpha), w)$. Moreover, $\operatorname{diag}\left[{ }^{t} \psi_{v}(\alpha),{ }^{t} \varphi_{v}(\alpha)\right]$, being the $v$-component of ${ }^{t} M(h(\alpha), w)$, represents $\iota_{w}(\alpha)$ on the $v$-th factor of $\left(\mathbf{C}^{l}\right)^{\mathbf{a}}$ (see §4.3). By Lemma 4.13, $\psi_{v}(a)$ and $\varphi_{v}(a)$ have algebraic entries for every $a \in Y$. Take $B_{v} \in G L_{m_{v}}(\overline{\mathbf{Q}})$ and $C_{v} \in G L_{n_{v}}(\overline{\mathbf{Q}})$ so that

$$
\begin{align*}
& B_{v} \psi_{v}(a) B_{v}^{-1}=\operatorname{diag}\left[\psi_{v 1}(a), \ldots, \psi_{v m_{v}}(a)\right]  \tag{11.3a}\\
& C_{v} \varphi_{v}(a) C_{v}^{-1}=\operatorname{diag}\left[\varphi_{v 1}(a), \ldots, \varphi_{v n_{v}}(a)\right] \tag{11.3b}
\end{align*}
$$

for every $a \in Y$ with $\psi_{v i}, \varphi_{v j} \in J_{Y}$. In Case SP we have $\psi_{v}=\varphi_{v}$ and $m_{v}=n_{v}=$ $n$, and so we take $\psi_{v i}=\varphi_{v i}$ and $B_{v}=C_{v}$.

We now define an element $\mathfrak{p}(w)=\left(\mathfrak{p}_{v}(w)\right)_{v \in \mathbf{b}}$ of $\prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C})$ in Cases UT and UB by

$$
\begin{align*}
\mathfrak{p}_{v \rho}(w) & =B_{v}^{-1} \operatorname{diag}\left[p_{Y}\left(\psi_{v 1}, \Phi\right), \ldots, p_{Y}\left(\psi_{v m_{v}}, \Phi\right)\right] B_{v}  \tag{11.4a}\\
\mathfrak{p}_{v}(w) & =C_{v}^{-1} \operatorname{diag}\left[p_{Y}\left(\varphi_{v 1}, \Phi\right), \ldots, p_{Y}\left(\varphi_{v n_{v}}, \Phi\right)\right] C_{v} \tag{11.4b}
\end{align*}
$$

where $v \in \mathbf{a}$. In Case SP we define $\mathfrak{p}_{v}(w)$ for each $v \in \mathbf{a}$ by (11.4b), and so $\mathfrak{p}(w)$ is an element of $G L_{n}(\mathbf{C})^{\mathbf{a}}$. Here recall the convention made in $\S \S 3.5$ and 5.1 that $\mathbf{b}=\mathbf{a} \rho \cup \mathbf{a}=J_{K}$ in Cases UT and UB, and $\mathbf{a}$ is identified with a fixed CM-type of $K$; also $n_{v \rho}=m_{v}$ for $v \in \mathbf{a}$. As we said in $\S 3.3$, we put $G L_{0}(\mathbf{C})=\{1\}$. Thus we define $\mathfrak{p}_{v}(w)$ to be the element 1 of $G L_{0}(\mathbf{C})$ if $n_{v}=0$. If we replace $C_{v}$ by another matrix $C_{v}^{\prime}$ in (11.3a), then $C_{v}^{\prime} C_{v}^{-1}$ is diagonal, and hence $\mathfrak{p}_{v}(w)$ does not change. Thus $\mathfrak{p}(w)$ is independent of the choice of $B_{v}$ and $C_{v}$. Also, if we view $p_{Y}($,$) as a$ coset of $\mathbf{C}^{\times} / \overline{\mathbf{Q}}^{\times}$, then we view $\mathfrak{p}(w)$ as a coset of $\prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C}) / \prod_{v \in \mathbf{b}} G L_{n_{v}}(\overline{\mathbf{Q}})$.

It may happen that the same CM-point can be obtained from different $(Y, h)$. However, we have:
11.5. Proposition. (1) The coset $\mathfrak{p}_{v}(w) G L_{n_{v}}(\overline{\mathbf{Q}})$ is determined by $w$ independently of the choice of $(Y, h)$.
(2) $\mathfrak{p}(\beta w) \mathcal{G}=M_{\beta}(w) \mathfrak{p}(w) \mathcal{G}$ for every $\beta \in \widetilde{G}_{+}$, where $\mathcal{G}=\prod_{v \in \mathbf{b}} G L_{n_{v}}(\overline{\mathbf{Q}})$.

Proof. Here we prove only (2). The proof of (1) will be completed in $\S \S 11.9$ and 11.10. Given $Y, h$, and $w$ as above, put $h^{\prime}(a)=\beta h(a) \beta^{-1}$. Then $\beta w$ is the fixed point of $h^{\prime}\left(Y^{u}\right)$ as observed in $\S 4.11$. Put $\xi_{v}=\lambda_{v}(\beta, w)$. Then we have $\lambda_{v}\left(h^{\prime}(\alpha), \beta w\right) \xi_{v}=\xi_{v} \lambda_{v}(h(\alpha), w)$ for every $\alpha \in Y^{u}$. Therefore, if we define $\psi_{v}^{\prime}$ and $\varphi_{v}^{\prime}$ by (4.37) with $h^{\prime}$ in place of $h$, then $B_{v} \xi_{v}^{-1} \psi_{v}^{\prime}(a) \xi_{v} B_{v}^{-1}=$ $\operatorname{diag}\left[\psi_{v 1}(a), \ldots, \psi_{v m_{v}}(a)\right]$ for every $a \in Y$. Thus $\mathfrak{p}_{v \rho}(\beta w)=\xi_{v} \mathfrak{p}_{v \rho}(w) \xi_{v}^{-1}$. A similar formula holds for $\mathfrak{p}_{v}(w)$ and $\mathfrak{p}_{v}(\beta w)$. This proves (2).
11.6. Putting $2 d=r[K: \mathbf{Q}]$, we are going to define an embedding of $\mathcal{H}$ into $\mathfrak{H}_{d}$. In Cases SP and UT this was done in $\S \S 6.2$ and 6.5 . Here we treat all cases uniformly. (Recall that $r=2 n$ in Cases SP and UT; also the present $d$ is the same as that of (4.9).) Since $(x, y) \mapsto \operatorname{Tr}_{K / \mathbf{Q}}\left(x \mathcal{T} y^{*}\right)$ with $\mathcal{T}$ of (4.13) is a nondegenerate alternating form, we can find a $\mathbf{Q}$-linear map $g: K_{r}^{1} \rightarrow \mathbf{Q}_{2 d}^{1}$ so that

$$
\begin{equation*}
\operatorname{Tr}_{K / \mathbf{Q}}\left(x \mathcal{T} y^{*}\right)=g(x) \eta_{d} \cdot{ }^{t} g(y) \quad\left(x, y \in K_{r}^{1}\right) \tag{11.5}
\end{equation*}
$$

Let $\left\{e_{k}\right\}_{k=1}^{2 d}$ be the standard basis of $\mathbf{Q}_{2 d}^{1}$. Given $z \in \mathcal{H}$, we consider the map $p_{z}$ : $\left(K_{\mathbf{a}}\right)_{r}^{1} \rightarrow \mathbf{C}^{d}$ of (4.23) and put $x_{k}=p_{z}\left(g^{-1}\left(e_{k}\right)\right)$. Now we can take ( $\mathbf{Q}, p_{z} \circ g^{-1}, \eta_{d}$ ) as $(W, q, \mathcal{T})$ in $\S \S 4.3$ and 4.4 , and we can write (4.18) and (4.21) in the form

$$
X\left(p_{z} \circ g^{-1}\right)=\left[\begin{array}{lll}
x_{1} & \cdots & x_{2 d}  \tag{11.6}\\
\bar{x}_{1} & \cdots & \bar{x}_{2 d}
\end{array}\right]=\left[\begin{array}{cc}
\kappa & 0 \\
0 & \bar{\kappa}
\end{array}\right]\left[\begin{array}{cc}
Z & 1_{d} \\
\bar{Z} & 1_{d}
\end{array}\right]
$$

with $\kappa \in G L_{d}(\mathbf{C})$ and $Z \in \mathfrak{H}_{d}$. Put $\kappa=\kappa(z)$ and $Z=\varepsilon(z)$. Then

$$
\varepsilon(z)=\kappa(z)^{-1}\left[\begin{array}{lll}
x_{1} & \cdots & x_{d}
\end{array}\right], \quad \kappa(z)=\left[\begin{array}{lll}
x_{d+1} & \cdots & x_{2 d} \tag{11.7}
\end{array}\right] .
$$

Since $p_{z}$ is holomorphic in $z$ as noted in $\S 4.7$, we see that both $\kappa(z)$ and $\varepsilon(z)$ are holomorphic in $z$. Thus we have a holomorphic embedding

$$
\begin{equation*}
\varepsilon: \mathcal{H} \rightarrow \mathfrak{H}_{d} . \tag{11.8}
\end{equation*}
$$

Now $\left(p_{z} \circ g^{-1}\right)\left(\sum_{k=1}^{m} c_{k} e_{k}\right)=\kappa(z)\left[\varepsilon(z) \quad 1_{d}\right] c$ for every $c \in \mathbf{Q}^{2 d}$, that is,

$$
\begin{equation*}
p_{z}(x)=\kappa(z)\left[\varepsilon(z) \quad 1_{d}\right] \cdot{ }^{t} g(x) \quad\left(x \in K_{r}^{1}\right) \tag{11.9}
\end{equation*}
$$

For $\alpha \in K_{r}^{r}$ define $\widetilde{\alpha} \in \mathbf{Q}_{2 d}^{2 d}$ by $g(x \alpha)=g(x) \widetilde{\alpha}$. From (11.5) we see that $\widetilde{\alpha} \in S p(d, \mathbf{Q})$ if $\alpha \in U(\mathcal{T})=G$. Also, from (4.31) and (11.9) we obtain

$$
\left.\begin{array}{l}
{ }^{t} M(\alpha, z) \kappa(\alpha z)[\varepsilon(\alpha z)  \tag{11.10}\\
=\kappa(z)\left[\begin{array} { l l } 
{ } \\
{ = }
\end{array} \left(x(x)={ }^{t} M(\alpha, z) p_{\alpha z}(x)=p_{z}(x \alpha)\right.\right. \\
=\kappa(z) \cdot t(x)=\kappa(z) \cdot{ }^{t} \mu(\widetilde{\alpha}, \varepsilon(z))[\widetilde{\alpha}(\varepsilon(z)) \\
1
\end{array}\right] \cdot{ }^{t} g(x) .
$$

This proves that $\varepsilon(\alpha z)=\widetilde{\alpha}(\varepsilon(z))$ and

$$
\begin{equation*}
\mu(\widetilde{\alpha}, \varepsilon(z))={ }^{t} \kappa(\alpha z) M(\alpha, z) \cdot{ }^{t} \kappa(z)^{-1} \quad(\alpha \in U(\mathcal{T}), z \in \mathcal{H}) \tag{11.11}
\end{equation*}
$$

To avoid confusion, we hereafter denote by $\mathfrak{A}_{0}(\overline{\mathbf{Q}})$ (resp. $\mathfrak{A}_{s}(\overline{\mathbf{Q}})$ ) the field $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ (resp. the vector space $\mathcal{A}_{s}(\overline{\mathbf{Q}})$ ) defined on $\mathfrak{H}_{d}$ with respect to $S p(d, \mathbf{Q})$, where $s$ is a $\mathbf{Q}$-rational representation of $G L_{d}(\mathbf{C})$. Then we note a simple fact:
(11.12) If $g \in \mathfrak{A}_{0}(\overline{\mathbf{Q}})$ and $g \circ \varepsilon$ is finite, then $g \circ \varepsilon \in \mathcal{A}_{0}(\overline{\mathbf{Q}})$.

To prove this, take a CM-point $w$ on $\mathcal{H}$ fixed by $h\left(Y^{u}\right)$ with $h: Y \rightarrow K_{r}^{r}$ satisfying (4.35). Define $\widetilde{h}: Y \rightarrow \mathbf{Q}_{2 d}^{2 d}$ by $\widetilde{h}(a)=\widetilde{h(a)}$. Then (4.35) together with (11.5) shows that $\widetilde{h}\left(a^{\rho}\right)=\eta_{d} \cdot t \tilde{h}(a) \eta_{d}^{-1}$, and $\varepsilon(w)$ is the fixed point of $\widetilde{h}\left(Y^{u}\right)$. Thus $\varepsilon(w)$ is a CM-point on $\mathfrak{H}_{d}$. Therefore (11.12) follows from Proposition 11.2 (2).

We note here a basic fact on 1-forms on an abelian variety:
11.7. Lemma. Let $A$ be an abelian variety defined over a subfield $k$ of $\mathbf{C}$ which is algebraic over the field of moduli of $(A, \mathcal{C})$ with a polarization $\mathcal{C}$ of $A$. Suppose that $A$ is isomorphic to $\mathbf{C}^{d} /\left(\left[\begin{array}{ll}Z_{0} & 1_{d}\end{array}\right] \mathbf{Z}^{2 d}\right)$ with $Z_{0} \in \mathfrak{H}_{d}$. Then there exists a $(d \times d)$ matrix $P$ whose columns belong to $\mathfrak{A}_{\sigma}(\overline{\mathbf{Q}})$ and such that $P$ is finite and invertible at $Z_{0}$, where $\sigma$ is the identity map of $G L_{d}(\mathbf{C})$ onto itself. Moreover, with any such $P$ define 1-forms $\xi_{1}, \ldots, \xi_{n}$ on $\mathbf{C}^{d} /\left(\left[\begin{array}{ll}Z_{0} & 1_{d}\end{array}\right] \mathbf{Z}^{2 d}\right)$ by

$$
\left[\begin{array}{c}
\xi_{1} \\
\vdots \\
\xi_{d}
\end{array}\right]=\pi \cdot{ }^{t} P\left(Z_{0}\right)\left[\begin{array}{c}
d u_{1} \\
\vdots \\
d u_{d}
\end{array}\right]
$$

where $u_{1}, \ldots, u_{d}$ are the coordinate functions on $\mathbf{C}^{d}$. Then $\xi_{1}, \ldots, \xi_{d}$, viewed as 1 -forms on $A$, form a basis of holomorphic 1-forms rational over the algebraic closure $\bar{k}$ of $k$.

This is a simplified version of [S98, Theorem 30.3]. The existence of $P$ is a special case of Proposition 9.11. Notice that $\bar{k}$ does not depend on the choice of $\mathcal{C}$.
11.8. We now restrict our discussion to Case UB. Take $\sigma$ and $P$ as in Lemma 11.7 with any fixed point $Z_{0}$ on $\mathfrak{H}_{d}$ in the setting of $\S 11.6$. By Lemma 4.12, we can take $b \in K$ so that $K=\mathbf{Q}(b)$ and $b b^{\rho}=1$; put $\alpha=b 1_{r}$ and

$$
U(z)=V(\varepsilon(z)), \quad V(Z)=P(Z)^{-1}\left(P \|_{\sigma} \widetilde{\alpha}\right)(Z) \quad\left(z \in \mathcal{H}, Z \in \mathfrak{H}_{d}\right)
$$

By Proposition 9.11 (4) the entries of $V$ belong to $\mathfrak{A}_{0}(\overline{\mathbf{Q}})$, and hence the entries of $U$ belong to $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ by (11.12). From (3.17), (3.37), (4.9), (4.10), and (4.29) we see that $M(\alpha, z)=\Psi(b)$, which combined with (11.11) shows that $\mu(\widetilde{\alpha}, \varepsilon(z))={ }^{t} \kappa(z) \Psi(b)$. ${ }^{t} \kappa(z)^{-1}$, and hence $U(z)^{-1}=X(z)^{-1} \Psi(b) X(z)$ with $X(z)={ }^{t} \kappa(z)^{-1} P(\varepsilon(z))$. Since $K=\mathbf{Q}[b]$, we see that $a \mapsto X(z)^{-1} \Psi(a) X(z)$ is a ring-injection of $K$ into $\mathcal{A}_{0}(\overline{\mathbf{Q}})_{d}^{d}$. Therefore we can find an element $W$ of $G L_{d}\left(\mathcal{A}_{0}(\overline{\mathbf{Q}})\right)$ such that $X(z)^{-1} \Psi(a) X(z)=$ $W^{-1} \Psi(a) W$ for every $a \in K$. In view of (4.10) we have

$$
\begin{equation*}
{ }^{t} \kappa(z)^{-1} P(\varepsilon(z))=X(z)=\operatorname{diag}\left[S_{v}, R_{v}\right]_{v \in \mathbf{a}} \cdot W(z) \tag{11.13}
\end{equation*}
$$

with square matrices $S_{v}$ and $R_{v}$ of size $m_{v}$ and $n_{v}$, whose entries are meromorphic functions on $\mathcal{H}$. Employing (11.11) and (4.29), we easily see that

$$
\begin{equation*}
S_{v}(\gamma(z))=\lambda_{v}(\gamma, z) S_{v}(z), \quad R_{v}(\gamma(z))=\mu_{v}(\gamma, z) R_{v}(z) \quad(v \in \mathbf{a}) \tag{11.14}
\end{equation*}
$$

if $\gamma$ belongs to a sufficiently small congruence subgroup $\Gamma$ of $G$.
11.9. We took $\mathcal{F}(\Omega)$ in $\S 11.1$ with a PEL-type $\Omega=\{K, \Psi, L, \mathcal{T}\}$. We now assume that $L=g^{-1}\left(\mathbf{Z}_{2 d}^{1}\right)$ with $g: K_{r}^{1} \rightarrow \mathbf{Q}_{2 d}^{1}$ as in $\S 11.6$. Let $w \in \mathcal{H}_{\mathrm{CM}}$. From Lemma 4.13 we see that $p_{w}(x)$ has algebraic components for every $x \in K_{r}^{1}$, and hence (11.6) shows that the entries of $\varepsilon(w)$ and $\kappa(w)_{v}$ for every $v \in \mathbf{a}$ are all algebraic.

Let $\mathcal{W}_{0}$ be the set of all CM-points on $\mathcal{H}$ where both $P \circ \varepsilon$ and $W$ are finite and invertible. Then, by (11.13), $R_{v}$ and $S_{v}$ are finite and invertible at every point of $\mathcal{W}_{0}$.

To prove (1) of Proposition 11.5 in Case UB, we first assume that $w \in \mathcal{W}_{0}$. Take ( $A_{w}, \iota^{\prime}$ ) of type ( $Y, \Phi$ ) we considered in $\S 11.3$. Recall that $A_{w}$ is isomorphic to $\mathbf{C}^{d} / p_{w}(L)$. Also, we see that $u \mapsto \kappa(w)^{-1} u$ for $u \in \mathbf{C}^{d}$ gives an isomorphism of $\mathbf{C}^{d} / p_{w}(L)$ onto $\mathbf{C}^{d} /\left[\varepsilon(w) \quad 1_{d}\right] \mathbf{Z}^{2 d}$. By Lemma 11.7 we obtain $\overline{\mathbf{Q}}$-rational 1-forms $\xi_{k}$ on $A_{w}$ by putting

$$
\left[\begin{array}{c}
\xi_{1} \\
\vdots \\
\xi_{d}
\end{array}\right]=\pi \cdot{ }^{t} P(\varepsilon(w)) \kappa(w)^{-1}\left[\begin{array}{c}
d u_{1} \\
\vdots \\
d u_{d}
\end{array}\right]
$$

where $u_{1}, \ldots, u_{d}$ are the standard coordinate functions on $\mathbf{C}^{d}$.
We have a decomposition $\mathbf{C}^{d}=\bigoplus_{v \in \mathbf{a}} V_{v}$ with $V_{v}$, isomorphic to $\mathbf{C}^{r}$, on which $\Psi_{v}$ of (24.1b) acts (see $\S 4.3$ ). Let $x_{1}^{v}, \ldots, x_{m_{v}}^{v}, y_{1}^{v}, \ldots, y_{n_{v}}^{v}$ be the coordinate
functions on $V_{v}$. Then these for all $v$ are renamings of the $u_{k}$. Since $W(w)$ is $\overline{\mathrm{Q}}$-rational, we find that the components of

$$
\pi \cdot{ }^{t} S_{v}(w)\left[\begin{array}{c}
d x_{1}^{v}  \tag{11.15}\\
\vdots \\
d x_{m_{v}}^{v}
\end{array}\right] \quad \text { and } \quad \pi \cdot{ }^{t} R_{v}(w)\left[\begin{array}{c}
d y_{1}^{v} \\
\vdots \\
d y_{n_{v}}^{v}
\end{array}\right]
$$

correspond to $\overline{\mathbf{Q}}$-rational 1-forms on $A_{w}$. Take $B_{v}$ as in (11.3a) and put

$$
\left[\begin{array}{c}
d u_{1}^{v} \\
\vdots \\
d u_{m_{v}}^{v}
\end{array}\right]={ }^{t} B_{v}^{-1}\left[\begin{array}{c}
d x_{1}^{v} \\
\vdots \\
d x_{m_{v}}^{v}
\end{array}\right] .
$$

Then we find that $\delta \iota(a) d u_{i}^{v}=\psi_{v i}(a) d u_{i}^{v}$. Now the periods of $d x_{i}^{v}$ are entries of $p_{w}(L)$, which are algebraic, and hence so are the periods of $d u_{i}^{v}$. Therefore by (11.1) we see that $\pi p_{Y}\left(\psi_{v i}, \Phi\right) d u_{i}^{v}$, viewed as a 1 -form on $A_{w}$, must be $\overline{\mathbf{Q}}$-rational. Comparing this result with (11.15), we find the first of the following two inclusions:

$$
\begin{aligned}
& \operatorname{diag}\left[p_{Y}\left(\psi_{v 1}, \Phi\right), \ldots, p_{Y}\left(\psi_{v m_{v}}, \Phi\right)\right] \in B_{v} S_{v}(w) G L_{m_{v}}(\overline{\mathbf{Q}}) \\
& \operatorname{diag}\left[p_{Y}\left(\varphi_{v 1}, \Phi\right), \ldots, p_{Y}\left(\varphi_{v n_{v}}, \Phi\right)\right] \in C_{v} R_{v}(w) G L_{n_{v}}(\overline{\mathbf{Q}})
\end{aligned}
$$

The latter can be shown similarly. These can be written

$$
\begin{equation*}
\mathfrak{p}_{v \rho}(w) \in S_{v}(w) G L_{m_{v}}(\overline{\mathbf{Q}}) \quad \text { and } \quad \mathfrak{p}_{v}(w) \in R_{v}(w) G L_{n_{v}}(\overline{\mathbf{Q}}) \tag{11.16}
\end{equation*}
$$

for every $v \in \mathbf{a}$. Now $S_{v}(w)$ and $R_{v}(w)$ depend only on $w$ (and $P$ ), and are independent of the choice of ( $Y, h$ ). Therefore we obtain (1) of Proposition 11.5 in Case UB at least for $w \in \mathcal{W}_{0}$.

Now take any CM-point $w$. Then we can find $\beta \in \widetilde{G}_{+}$so that $\beta^{-1} w \in \mathcal{W}_{0}$. Put $w_{1}=\beta^{-1} w$. By (2) of Proposition 11.5, $\mathfrak{p}(w) \mathcal{G}=M_{\beta}\left(w_{1}\right) \mathfrak{p}\left(w_{1}\right) \mathcal{G}$, where $\mathfrak{p}(w)$ and $\mathfrak{p}\left(w_{1}\right)$ are defined with a fixed $(Y, h)$ and $\left(Y, h^{\prime}\right)$ as in the proof there. Now $\mathfrak{p}\left(w_{1}\right)$ is independent of the choice of $(Y, h)$, and therefore the same is true for $\mathfrak{p}(w)$. This completes the proof of Proposition 11.5 in Case UB.
11.10. Let us now prove (1) of Proposition 11.5 in Cases SP and UT. In Case SP, (6.12) shows that (11.9) is true for $z \in \mathfrak{H}_{n}^{\mathrm{a}}$ with $\kappa(z)=B^{-1}$, and ( 6.11 b ) is exactly (11.11). Now we consider $R_{v}$ of Proposition 9.11 belonging to $\mathcal{A}_{\sigma_{v}}(\overline{\mathbf{Q}})$. Let $\mathfrak{R}(z)=$ $\operatorname{diag}\left[R_{v}(z)\right]_{v \in \mathbf{a}}$ and $W(z)=\mathfrak{R}(z)^{-1} \cdot{ }^{t} B P(\varepsilon(z))$, where $P$ is the function on $\mathfrak{H}_{d}$ as in Lemma 11.7. Then $W$ has entries in $\mathcal{A}_{0}$. Checking the Fourier coefficients, we see that the entries of $W$ are $\overline{\mathbf{Q}}$-rational. Since ${ }^{t} B P(\varepsilon(z))=\operatorname{diag}\left[R_{v}(z)\right]_{v \in \mathbf{a}} W(z)$, (11.13) is true with the present symbols, if we ignore $S_{v}$. Therefore we can repeat our argument of $\S 11.9$ in Case SP to find that (11.16) is true in the present setting, and we obtain the desired fact.

In Case UT, combining (6.24b) with (6.12), we obtain an element $E$ of $G L_{d}(\overline{\mathbf{Q}})$ and an embedding $\varepsilon$ of $\mathcal{H}=\mathcal{H}_{n}^{\text {a }}$ into $\mathfrak{H}_{d}$ such that $p_{z}(x)=E\left[\varepsilon(z) 1_{d}\right] \cdot{ }^{t} g(x)$ with a map $g: K_{2 n}^{1} \rightarrow \mathbf{Q}_{2 d}^{1}$ satisfying (11.5) with $\mathcal{T}=\eta_{n}$. (Using the symbols of $\S \S 6.2$ and $6.5, g^{-1} \circ h^{-1}$ gives the present $g$.) Therefore, considering the functions $R_{v}$ of Proposition 9.11, we can handle Case UT in the same manner as in Case SP.

As a by-product of this reasoning we obtain:
11.11. Proposition. Let $\omega$ be a $\overline{\mathbf{Q}}$-rational representation of $G L_{n}(\mathbf{C})^{\mathbf{b}}$ in Cases $S P$ and UT. Let $\mathcal{W}$ be a subset of $\mathcal{H}_{\mathrm{CM}}$ dense in $\mathcal{H}$. Then an element $f$
of $\mathcal{A}_{\omega}$ belongs to $\mathcal{A}_{\omega}(\overline{\mathbf{Q}})$ if and only if $\omega\left(\mathfrak{p}(w)^{-1}\right) f(w)$ is $\overline{\mathbf{Q}}$-rational for every $w \in \mathcal{W}$ where $f$ is finite.

Proof. Take $R=\left(R_{v}\right)_{v \in \mathbf{b}}$ with $R_{v} \in \mathcal{A}_{\sigma_{v}}(\overline{\mathbf{Q}})$ as in Proposition 9.11. Let $\mathcal{W}^{\prime}$ be the subset of $\mathcal{W}$ consisting of the points $w$ where $R_{v}$ is finite and invertible for every $v \in \mathbf{b}$. Then $\mathcal{W}^{\prime}$ is dense in $\mathcal{H}$. Put $g=\omega(R)^{-1} f$. Then $g$ has components in $\mathcal{A}_{0}$. Suppose $\omega\left(\mathfrak{p}(w)^{-1}\right) f(w)$ is $\overline{\mathbf{Q}}$-rational for every $w \in \mathcal{W}$ where $f$ is finite. Then $g(w)=\omega\left(R(w)^{-1} \mathfrak{p}(w)\right) \omega\left(\mathfrak{p}(w)^{-1}\right) f(w)$. The reasoning in $\S 11.10$ shows that (11.16) is true in Cases SP and UT with the present $R_{v}$ and $S_{v}=R_{v \rho}$, that is, $R(w)^{-1} \mathfrak{p}(w)$ is $\overline{\mathbf{Q}}$-rational for every $w \in \mathcal{W}^{\prime}$, and hence $g(w)$ is $\overline{\mathbf{Q}}$-rational for every $w \in \mathcal{W}^{\prime}$ where $g$ is finite. By Proposition 11.2 (2), $g$ has components in $\mathcal{A}_{0}(\overline{\mathbf{Q}})$. Since $f=\omega(R) g$, we see that $f \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$. This proves the 'if'-part. Conversely, if $f \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$, then $g \in \mathcal{A}_{0}(\overline{\mathbf{Q}})$. For every $w \in \mathcal{W}$ we can choose $R$ so that $R_{v}$ is finite and invertible at $w$ for every $v \in \mathbf{b}$. If $f$ is finite at $w$, then so is $g$, and $g(w)$ is $\overline{\mathbf{Q}}$-rational. We have $\omega\left(\mathfrak{p}(w)^{-1}\right) f(w)=\omega\left(\mathfrak{p}(w)^{-1} R(w)\right) g(w)$, which is $\overline{\mathbf{Q}}$-rational. This completes the proof.
11.12. Take a $\overline{\mathbf{Q}}$-rational representation $\omega: \prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C}) \rightarrow G L(V)$ with a finite-dimensional complex vector space $V$ with a $\overline{\mathbf{Q}}$-rational structure. Then for each CM-point $w$ we put

$$
\begin{align*}
& \mathfrak{P}_{\omega}(w)=\omega(\mathfrak{p}(w)),  \tag{11.17}\\
& \mathfrak{P}_{k}(w)=\operatorname{det}(\mathfrak{p}(w))^{k} \quad\left(k \in \mathbf{Z}^{\mathbf{b}}\right), \tag{11.17a}
\end{align*}
$$

and denote by $\mathcal{A}_{\omega}(\overline{\mathbf{Q}})$ the set of all $f \in \mathcal{A}_{\omega}$ such that $\mathfrak{P}_{\omega}(w)^{-1} f(w)$ is $\overline{\mathbf{Q}}$-rational for every $w \in \mathcal{H}_{\mathrm{CM}}$ where $f$ is finite. We call such an $f \overline{\mathbf{Q}}$-rational, and put $\mathcal{M}_{\omega}(\overline{\mathbf{Q}})=\mathcal{M}_{\omega} \cap \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$. Proposition 11.11 shows that this is consistent with what we already have in Cases SP and UT. Thus the point of our definition in this subsection is mainly in Case UB.
11.13. Proposition (Case UB). (1) If $f \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$ and $\beta \in \widetilde{G}_{+}$, then both $\left.f\right|_{\omega} \beta$ and $f \|_{\omega} \beta$ belong to $\mathcal{A}_{\omega}(\overline{\mathbf{Q}})$.
(2) Let $\mathcal{W}$ be a subset of $\mathcal{H}_{\mathrm{CM}}$ dense in $\mathcal{H}$. Then an element $f$ of $\mathcal{A}_{\omega}$ belongs to $\mathcal{A}_{\omega}(\overline{\mathbf{Q}})$ if $\mathfrak{P}_{\omega}(w)^{-1} f(w)$ is $\overline{\mathbf{Q}}$-rational for every $w \in \mathcal{W}$ where $f$ is finite.

Proof. Assertion (1) follows immediately from our definition of $\mathcal{A}_{\omega}(\overline{\mathbf{Q}})$ and Proposition 11.5 (2). Assertion (2) will be proved in the proof of the following proposition.
11.14. Proposition (Case UB). For each fixed $v \in \mathbf{b}$ define two $G L_{n_{v}}(\mathbf{C})$ valued representations $\sigma_{v}$ and $\tau_{v}$ of $\prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C})$ by

$$
\begin{equation*}
\sigma_{v}(x)=x_{v} \quad \text { and } \quad \tau_{v}(x)=\operatorname{det}(x)^{\mathbf{b}} x_{v} \quad \text { for } \quad x \in \prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C}) \tag{11.18}
\end{equation*}
$$

where $a^{\mathbf{b}}=\prod_{v \in \mathbf{b}} a_{v}$. Given $z_{0} \in \mathcal{H}$, there exist two sets of matrix-valued functions $\left\{R_{v}\right\}_{v \in \mathbf{b}}$ and $\left\{Q_{v}\right\}_{v \in \mathbf{b}}$ on $\mathcal{H}$ with the following properties:
(1) $R_{v}$ and $Q_{v}$ are square matrices of size $n_{v}$, and the columns of $R_{v}$ (resp. $Q_{v}$ ) belong to $\mathcal{A}_{\sigma_{v}}(\overline{\mathbf{Q}})\left(\mathrm{resp} . \mathcal{M}_{\tau_{v}}(\overline{\mathbf{Q}})\right.$ ).
(2) $R_{v}$ is finite at $z_{0}$, and $\operatorname{det}\left(R_{v}\left(z_{0}\right)\right) \operatorname{det}\left(Q_{v}\left(z_{0}\right)\right) \neq 0$ for every $v \in \mathbf{b}$.

If $n_{v}=0$, we have $G L_{n_{v}}(\mathbf{C})=\{1\}$ by our stipulation. Therefore, if $n_{v}=0$, we simply take $R_{v}$ and $Q_{v}$ to be the constant on $\mathcal{H}$ whose values are 1 in this trivial
group. Notice also that Proposition 9.11 gives similar results in Cases SP and UT in stronger forms.

Proof. Let the notation be as in $\S \S 11.8$ and 11.9. Put $R_{v \rho}=S_{v}$ for each $v \in \mathbf{a}$. Then (11.14) shows that $R_{v} \in \mathcal{A}_{\sigma_{v}}$ for every $v \in \mathbf{b}$. Moreover, by (11.16), $\mathfrak{p}_{v}(w)^{-1} R_{v}(w)$ is $\overline{\mathbf{Q}}$-rational for every $w \in \mathcal{W}_{0}$. Let $w_{1}$ be a CM-point where $R_{v}$ (for a fixed $v$ ) is finite. Take $\alpha \in G$ so that $w_{1}=\alpha w_{0}$ with $w_{0} \in \mathcal{W}_{0}$. Put $R_{v}^{\prime}=$ $R_{v} \|_{\sigma_{v}} \alpha$. Then $R_{v}^{-1} R_{v}^{\prime}$ has entries in $\mathcal{A}_{0}$. Take a CM-point $w$ such that both $w$ and $\alpha w$ belong to $\mathcal{W}_{0}$. Then $\left(R_{v}^{-1} R_{v}^{\prime}\right)(w) \in G L_{n_{v}}(\overline{\mathbf{Q}}) \mathfrak{p}_{v}(w)^{-1} \mu_{v}(\alpha, w)^{-1} R_{v}(\alpha w)=$ $G L_{n_{v}}(\overline{\mathbf{Q}}) \mathfrak{p}_{v}(\alpha w)^{-1} R_{v}(\alpha w)=G L_{n_{v}}(\overline{\mathbf{Q}})$. By Proposition $11.2(2), R_{v}^{-1} R_{v}^{\prime}$ has entries in $\mathcal{A}_{0}(\overline{\mathbf{Q}})$. Therefore $\mathfrak{p}_{v}\left(w_{1}\right)^{-1} R_{v}\left(w_{1}\right) \in G L_{n_{v}}(\overline{\mathbf{Q}}) \mathfrak{p}_{v}\left(w_{0}\right)^{-1} \mu_{v}\left(\alpha, w_{0}\right)^{-1} R_{v}\left(\alpha w_{0}\right) \subset$ $G L_{n_{v}}(\overline{\mathbf{Q}})\left(R_{v}^{-1} R_{v}^{\prime}\right)\left(w_{0}\right) \subset G L_{n_{v}}(\overline{\mathbf{Q}})$. This shows that the columns of $R_{v}$ belong to $\mathcal{A}_{\sigma_{v}}(\overline{\mathbf{Q}})$. Now given $z_{0} \in \mathcal{H}$, we can find $\beta \in G$ such that $R_{v}$ is finite and invertible at $\beta z_{0}$ for every $v \in \mathbf{b}$. Changing $R_{v}$ for $R_{v} \| \beta$, we obtain the desired elements of $\mathcal{A}_{\sigma_{v}}(\overline{\mathbf{Q}})$ in view of Proposition 11.13 (1).

Let us now prove (2) of Proposition 11.13. Given $\mathcal{W}$ and $f$ in that assertion, take $R=\left(R_{v}\right)_{v \in \mathbf{b}}$ as in the present proposition and put $g=\omega(R)^{-1} f$. Then we can show that $g$ has entries in $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ by the same technique as in the proof of Proposition 11.11. Now let $w$ be a CM-point where $f$ is finite. Choose $R$ so that $R_{v}$ is finite and invertible at $w$ for every $v \in \mathbf{b}$. Then $g$ is finite at $w$ and $g(w)$ is algebraic. Since $f(w)=\omega(R(w)) g(w)$, we see that $\omega(\mathfrak{p}(w))^{-1} f(w)$ is algebraic, so that $f \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$ as desired.

To prove the existence of $\left(Q_{v}\right)_{v \in \mathbf{b}}$, define $\tau: G L_{d}(\mathbf{C}) \rightarrow G L_{d}(\mathbf{C})$ by $\tau(x)=$ $\operatorname{det}(x) x$. By Proposition 9.11, given any $z_{0} \in \mathcal{H}$, there exists a $d \times d$-matrix $E$ whose columns belong to $\mathcal{M}_{\tau}(\mathbf{Q})$ and such that $\operatorname{det} E\left(\varepsilon\left(z_{0}\right)\right) \neq 0$. Put $S(z)=$ $\operatorname{det}(\kappa(z))^{-1} \cdot{ }^{t} \kappa(z)^{-1} E(\varepsilon(z))$ for $z \in \mathcal{H}$. From (11.11) and (4.29) we easily see that $S(\gamma(z))=j_{\gamma}(z)^{\mathbf{b}} M(\gamma, z) S(z)$ for every $\gamma$ in a congruence subgroup $\Gamma$ of $G$. Let $S_{v}$ be the submatrix of $S$ composed of the $n_{v}$ rows of $S$ corresponding to the component $\mu_{v}(\gamma, z)$ of $M(\gamma, z)$. Then $S_{v}(\gamma(z))=j_{\gamma}(z)^{\mathbf{b}} \mu_{v}(\gamma, z) S_{v}(z)$ for every $\gamma \in \Gamma$. Since $\operatorname{det} E\left(\varepsilon\left(z_{0}\right)\right) \neq 0$, we can find suitable $n_{v}$ columns of $S_{v}$ which make nonzero determinant at $z_{0}$. Call $Q_{v}$ the ( $n_{v} \times n_{v}$ )-matrix composed of those $n_{v}$ columns, which clearly belong to $\mathcal{A}_{\tau_{v}}$. Now by (11.13) and (11.16), for every $w \in \mathcal{W}_{0}$ we have

$$
P(\varepsilon(w)) \in^{t} \kappa(w) \operatorname{diag}\left[\mathfrak{p}_{v}(w)\right]_{v \in \mathbf{b}} G L_{d}(\overline{\mathbf{Q}})
$$

Since $P \in \mathcal{A}_{\sigma}(\overline{\mathbf{Q}})$, this together with Proposition 11.11 implies that

$$
\mathfrak{p}(\varepsilon(w)) \in^{t} \kappa(w) \operatorname{diag}\left[\mathfrak{p}_{v}(w)\right]_{v \in \mathbf{b}} G L_{d}(\overline{\mathbf{Q}}) .
$$

Put $\mathfrak{q}(w)=\prod_{v \in \mathbf{b}} \operatorname{det}\left(\mathfrak{p}_{v}(w)\right)$. Then $\operatorname{det}[\mathfrak{p}(\varepsilon(w))] \in \mathfrak{q}(w) \overline{\mathbf{Q}}$, and hence

$$
\mathfrak{q}(w)^{-1} \operatorname{diag}\left[\mathfrak{p}_{v}(w)\right]_{v \in \mathbf{b}}^{-1} S(w) \in G L_{d}(\overline{\mathbf{Q}}) \tau(\mathfrak{p}(\varepsilon(w)))^{-1} E(\varepsilon(w))=G L_{d}(\overline{\mathbf{Q}})
$$

Therefore $\tau_{v}(\mathfrak{p}(w))^{-1} Q_{v}(w)$ is algebraic for every $w \in \mathcal{W}_{0}$. By Proposition 11.13 (2), the columns of $Q_{v}$ belong to $\mathcal{M}_{\tau_{v}}(\overline{\mathbf{Q}})$. This complets the proof.
11.15. Proposition. For every $\omega$ as in $\S 11.12$ and every congruence subgroup $\Gamma$ of $G$ we have $\mathcal{M}_{\omega}(\Gamma)=\mathcal{M}_{\omega}(\Gamma, \overline{\mathbf{Q}}) \otimes \mathbf{C}$.

Proof. In Cases SP and UT we already proved a stronger result in Theorem 10.4. Thus our question here is in Case UB. We first prove that if $f_{1}, \ldots, f_{t}$ are
elements of $\mathcal{M}_{\omega}(\overline{\mathbf{Q}})$ linearly independent over $\overline{\mathbf{Q}}$, then they are linearly independent over $\mathbf{C}$. Indeed, for each $w \in \mathcal{H}_{\mathrm{CM}}$ put $X_{w}=\left\{x \in \mathbf{C}^{t} \mid \sum_{i=1}^{t} x_{i} \mathfrak{P}_{\omega}(w)^{-1} f_{i}(w)=0\right\}$; let $Y$ be the intersection of $X_{w}$ for all such $w$. Clearly we can omit $\mathfrak{P}_{\omega}(w)^{-1}$ in the definition of $X_{w}$. Since $\mathfrak{P}_{\omega}(w)^{-1} f_{i}(w)$ is algebraic, each $X_{w}$, as well as $Y$, is a vector subspace of $\mathbf{C}^{t}$ defined over $\overline{\mathbf{Q}}$. If $Y$ contains a nonzero element $c$, then $\sum_{i=1}^{t} c_{i} f_{i}(w)=0$ for all CM-points $w$ so that $\sum_{i=1}^{t} c_{i} f_{i}=0$. Since the $f_{i}$ are linearly independent over $\overline{\mathbf{Q}}, Y$ has no such $c$ in $\overline{\mathbf{Q}}^{t}$, so that $Y=\{0\}$, which shows that the $f_{i}$ are linearly independent over $\mathbf{C}$.

Our next task is to show that $\mathcal{M}_{\omega}$ can be spanned by $\mathcal{M}_{\omega}(\overline{\mathbf{Q}})$ over $\mathbf{C}$. Denote $\mathcal{A}_{\omega}$ and $\mathcal{M}_{\omega}$ by $\mathcal{A}_{\kappa \mathbf{b}}$ and $\mathcal{M}_{\kappa \mathbf{b}}$ if $\omega(x)=\operatorname{det}(x)^{\kappa \mathbf{b}}$ with an integer $\kappa$ (see (5.4b)). In this special case our assertion is

$$
\begin{equation*}
\mathcal{M}_{\kappa \mathbf{b}}=\mathcal{M}_{\kappa \mathbf{b}}(\overline{\mathbf{Q}}) \otimes \mathbf{C} \tag{11.19}
\end{equation*}
$$

We shall prove this after the proof of Proposition 14.8. Assuming this result, we are going to prove the desired result for $\mathcal{M}_{\omega}$. We may assume that $\omega$ is irreducible. Then there is an integer $e$ such that $\omega(c y)=c^{e} \omega(y)$ for $c \in \mathbf{C}^{\times}$. Take $Q=\left(Q_{v}\right)_{v \in \mathbf{b}}$ as in Proposition 11.14 with any $z_{0} \in \mathcal{H}$. Then $\omega(Q) \in \mathcal{M}_{\omega^{\prime}}(\overline{\mathbf{Q}})$ with $\omega^{\prime}(x)=$ $\operatorname{det}(x)^{e \mathbf{b}} \omega(x)$. Take a positive integer $p$ so that $p>e$ and $\operatorname{det}(x)^{p \mathbf{b}} \omega(x)^{-1}$ is a polynomial in $x$; put $\zeta(x)=\operatorname{det}(x)^{-p \mathbf{b}} \omega(x)$. Then $\zeta(Q) \in \mathcal{A}_{\xi}(\overline{\mathbf{Q}})$ with $\xi(x)=$ $\operatorname{det}(x)^{-\kappa \mathbf{b}} \omega(x)$, where $\kappa=p(1+r|\mathbf{a}|)-e$. Given $f \in \mathcal{M}_{\omega}$, put $g=\zeta(Q)^{-1} f$; then $g$ has components in $\mathcal{A}_{\kappa \mathbf{b}}$. Since $\zeta^{-1}$ is a polynomial function, $g$ is holomorphic, so that $g$ has components in $\mathcal{M}_{\kappa \mathbf{b}}$, or rather, $g \in\left(\mathcal{M}_{\kappa \mathbf{b}}\right)^{t}$, where $t$ is the dimension of the representation space of $\omega$. Take a $\overline{\mathbf{Q}}$-basis $\left\{c_{\nu}\right\}$ of $\mathbf{C}$. By (11.18) we can express $g$ as a finite sum $g=\sum_{\nu} c_{\nu} g_{\nu}$ with $g_{\nu} \in\left(\mathcal{M}_{\kappa \mathbf{b}}(\overline{\mathbf{Q}})\right)^{t}$. Then $f=\sum_{\nu} c_{\nu} f_{\nu}$ with $f_{\nu}=\zeta(Q) g_{\nu}$. Clearly $f_{\nu} \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$. Suppose $f=\sum_{\nu} c_{\nu} f_{\nu}^{\prime}$ with $f_{\nu}^{\prime} \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$. Let $w$ be a CM-point where $f_{\nu}$ and $f_{\nu}^{\prime}$ are finite. Then $0=\sum_{\nu} c_{\nu} \mathfrak{P}_{\omega}(w)^{-1}\left(f_{\nu}(w)-\right.$ $\left.f_{\nu}^{\prime}(w)\right)$, so that $f_{\nu}(w)=f_{\nu}^{\prime}(w)$. Since all such points $w$ form a dense subset of $\mathcal{H}$, we obtain $f_{\nu}=f_{\nu}^{\prime}$. This means that the $f_{\nu}$ are unique for $f$, once $\left\{c_{\nu}\right\}$ is chosen. Now, given $z_{0} \in \mathcal{H}$, take $Q$ so that $\operatorname{det} Q_{v}\left(z_{0}\right) \neq 0$ for every $v \in \mathbf{b}$. Then $\zeta(Q)$ is finite at $z_{0}$, so that $f_{\nu}$ is finite at $z_{0}$. Thus $f_{\nu}$ is finite everywhere, and hence $f_{\nu} \in \mathcal{M}_{\omega}(\overline{\mathbf{Q}})$. Now suppose $f \in \mathcal{M}_{\omega}(\Gamma)$. Then $f=f\left\|\gamma=\sum_{\nu} c_{\nu} f_{\nu}\right\| \gamma$ for every $\gamma \in \Gamma$. Since $f_{\nu} \| \gamma \in \mathcal{M}_{\omega}(\overline{\mathbf{Q}})$ by Proposition 11.13 (1), the uniqueness of $f_{\nu}$ just proved shows that $f_{\nu} \| \gamma=f_{\nu}$, and hence $f_{\nu} \in \mathcal{M}_{\omega}(\Gamma, \overline{\mathbf{Q}})$. This proves our proposition.
11.16. Before proceeding further, let us introduce some symbols. Let $L$ be a CM-field containing $K$. For $\alpha \in J_{K}$ and $\beta \in J_{L}$ we denote by $\operatorname{Inf}_{L / K}(\alpha)$ the sum of all the elements of $J_{L}$ which coincide with $\alpha$ on $K$, and by $\operatorname{Res}_{L / K}(\beta)$ the restriction of $\beta$ to $K$. We then extend these to additive maps

$$
\begin{equation*}
\operatorname{Inf}_{L / K}: I_{K} \rightarrow I_{L}, \quad \operatorname{Res}_{L / K}: I_{L} \rightarrow I_{K} \tag{11.20}
\end{equation*}
$$

We recall here three basic properties of the period symbol $p_{K}$ (see [S98, Theorem 32.5]):

$$
\begin{align*}
& p_{K}\left(\xi, \operatorname{Res}_{L / K}(\zeta)\right)=p_{L}\left(\operatorname{Inf}_{L / K}(\xi), \zeta\right) \quad \text { if } \quad \xi \in I_{K} \quad \text { and } \zeta \in I_{L}  \tag{11.21}\\
& p_{K}\left(\operatorname{Res}_{L / K}(\zeta), \xi\right)=p_{L}\left(\zeta, \operatorname{Inf}_{L / K}(\xi)\right) \text { if } \xi \in I_{K} \text { and } \zeta \in I_{L}  \tag{11.22}\\
& p_{K}(\xi \rho, \eta)=p_{K}(\xi, \eta \rho)=p_{K}(\xi, \eta)^{-1} \quad \text { for every } \xi, \eta \in I_{K} \tag{11.23}
\end{align*}
$$

Here $p_{K}(\cdots)$ and $p_{L}(\cdots)$ are elements of $\mathbf{C}^{\times} / \overline{\mathbf{Q}}^{\times}$; also in the following proposition we view $\operatorname{det}\left(\mathfrak{p}_{v}(w)\right)$ for each $v \in \mathbf{b}$ as an element of $\mathbf{C}^{\times} / \overline{\mathbf{Q}}^{\times}$.
11.17. Theorem (Case UB). Let $\tau=\left\{\tau_{v}\right\}_{v \in \mathbf{a}}$ be a CM-type as in §3.5, and $\Psi$ the representation of $K$ as in (4.9) and (4.10); identify $\Psi$ with the element $\sum_{v \in \mathbf{a}}\left(m_{v} \tau_{v} \rho+n_{v} \tau_{v}\right)$ of $I_{K}$. Then the following assertions hold:
(1) $\operatorname{det}\left(\mathfrak{p}_{v}(w)\right) / \operatorname{det}\left(\mathfrak{p}_{v \rho}(w)\right)=p_{K}\left(\tau_{v}, \Psi\right)$ for every $w \in \mathcal{H}_{\mathrm{CM}}$. In particular, if $m_{v}=n_{v}$ for every $v \in \mathbf{a}$, then $\operatorname{det}\left(\mathfrak{p}_{v}(w)\right)=\operatorname{det}\left(\mathfrak{p}_{v \rho}(w)\right)$ for every $v \in \mathbf{a}$.
(2) Let $\omega(x)=\zeta(x) \prod_{v \in \mathbf{a}} \operatorname{det}\left(x_{v}\right)^{a_{v}} \operatorname{det}\left(x_{v \rho}\right)^{b_{v}}$ with $a, b \in \mathbf{Z}^{\mathbf{a}}$ such that $a_{v}+$ $b_{v}=0$ whenever $m_{v} n_{v} \neq 0$. Then $\mathcal{A}_{\omega}=\mathcal{A}_{\zeta}$ and $\mathcal{A}_{\omega}(\overline{\mathbf{Q}})=\mathfrak{q} \cdot \mathcal{A}_{\zeta}(\overline{\mathbf{Q}})$, where

$$
\mathfrak{q}=p_{K}\left(\sum_{v \in \mathbf{a}} a_{v} \tau_{v}-\sum_{v \in \mathbf{a}, n_{v}=0}\left(a_{v}+b_{v}\right) \tau_{v}, \Psi\right) .
$$

In particular, $\mathcal{A}_{\omega}(\overline{\mathbf{Q}})=\mathcal{A}_{\zeta}(\overline{\mathbf{Q}})$ if $m_{v}=n_{v}$ for every $v \in \mathbf{a}$.
Remark. In view of (11.23), we can replace $\Psi$ in these assertions by $\sum_{v \in \mathbf{a}}\left(n_{v}-\right.$ $\left.m_{v}\right) \tau_{v}$. Since $\mathfrak{q}$ can often be transcendental, the last assertion shows that the $\overline{\mathbf{Q}}$ rationality of automorphic forms in Case UB is far more complex than that in Cases SP and UT.

Proof. If $(Y, \Phi)$ is as in $\S 11.4$, then we have CM-types $\left(K_{i}, \Phi_{i}\right)$ such that $Y=\bigoplus_{i=1}^{t} K_{i}$ and $\Phi=\sum_{i=1}^{t} \Phi_{i}$ as in $\S 4.11$. Using the symbols of (11.3a, b), for each fixed $v \in \mathbf{a}$ we observe that the $\psi_{v i} \rho$ and the $\varphi_{v j}$ are exactly the elements of $I_{Y}$ whose restrictions to $K$ coincide with $\tau_{v}$. Thus $\sum_{i=1}^{m_{v}} \psi_{v i} \rho+\sum_{j=1}^{n_{v}} \varphi_{v j}=$ $\sum_{i=1}^{t} \operatorname{Inf}_{K_{i} / K}\left(\tau_{v}\right)$. Also $\Psi=\sum_{i=1}^{t} \operatorname{Res}_{K_{i} / K}\left(\Phi_{i}\right)$. Therefore from (11.4a, b), (11.23), (11.2), and (11.21) we obtain

$$
\begin{aligned}
\operatorname{det}\left(\mathfrak{p}_{v}(w)\right) / \operatorname{det}\left(\mathfrak{p}_{v \rho}(w)\right) & =\prod_{i=1}^{t} p_{K_{i}}\left(\operatorname{Inf}_{K_{i} / K}\left(\tau_{v}\right), \Phi_{i}\right) \\
& =\prod_{i=1}^{t} p_{K}\left(\tau_{v}, \operatorname{Res}_{K_{i} / K}\left(\Phi_{i}\right)\right)=p_{K}\left(\tau_{v}, \Psi\right)
\end{aligned}
$$

This proves the first half of (1). Suppose $m_{v}=n_{v}$ for every $v \in \mathbf{a}$; then $\Psi=$ $\sum_{i=1}^{t} n_{v}\left(\tau_{v} \rho+\tau_{v}\right)$, and hence $p_{K}\left(\tau_{v}, \Psi\right)=1$ for every $v \in \mathbf{a}$ by (11.23). This proves the latter half of (1). Let the notation be as in (2). From (3.23) and (4.34) we see that $\mathcal{A}_{\omega}=\mathcal{A}_{\zeta}$. (We already made the same type of observation in §5.4.) Now $\omega(x)=\zeta(x) \prod_{v \in \mathbf{a}}\left[\operatorname{det}\left(x_{v}\right) / \operatorname{det}\left(x_{v \rho}\right)\right]^{a_{v}} \operatorname{det}\left(x_{v \rho}\right)^{a_{v}+b_{v}}$, and $a_{v}+b_{v} \neq 0$ only if $m_{v} n_{v}=0$. If $n_{v}=0$, then $\operatorname{det}\left(\mathfrak{p}_{v}(w)\right)=1$, so that $\operatorname{det}\left(\mathfrak{p}_{v \rho}(w)\right)=p_{K}\left(-\tau_{v}, \Psi\right)$ by (1); if $m_{v}=0$, then $n_{v} \neq 0$ and $\operatorname{det}\left(\mathfrak{p}_{v \rho}(w)\right)=1$. Therefore by $(1), \mathfrak{P}_{\omega}(w)=$ $\mathfrak{q} \mathfrak{P}_{\zeta}(w)$ with $\mathfrak{q}$ given as above, and $\mathfrak{q}=1$ if $m_{v}=n_{v}$ for every $v \in \mathbf{a}$. Therefore we obtain (2) from our definion of $\mathcal{A}_{\omega}(\overline{\mathbf{Q}})$.
11.18. Proposition. With $\widetilde{G}=G p(1, F)=G L_{2}(F)$ let $w$ be a CM-point on $\mathcal{H}=\mathfrak{H}_{1}^{\mathbf{a}}$ obtained from $\left(K,\left\{\tau_{v}\right\}_{v \in \mathbf{a}}\right)$ and $w_{0} \in K$ as in Proposition 4.14. Then for $k \in \mathbf{Z}^{\mathbf{a}}$ we have $\mathfrak{P}_{k}(w)=p_{K}\left(\sum_{v \in \mathbf{a}} k_{v} \tau_{v}, \sum_{v \in \mathbf{a}} \tau_{v}\right)$.

Proof. In Proposition 4.14 we have seen that $\tau_{v}=\varphi_{v}$ and $\Phi=\sum_{v \in \mathbf{a}} \tau_{v}$, so that $\mathfrak{p}_{v}(w)=p_{K}\left(\tau_{v}, \Phi\right)$. Therefore $\mathfrak{P}_{k}(w)=\mathfrak{p}(w)^{k}=p_{K}\left(\sum_{v \in \mathbf{a}} k_{v} \tau_{v}, \Phi\right)$.
11.19. Proposition. The coset $p_{Y}(\alpha, \beta)$ of (11.2) can be represented by a real number.

Proof. It is sufficient to consider the case where $Y$ is a CM-field $K$. Let $K, \Phi, w_{0}$, and $w$ be as in Proposition 11.18; fix $u \in \mathbf{a}$ and put $\Phi^{\prime}=\Phi-\tau_{u}+\tau_{u} \rho$.

Then $\Phi^{\prime}$ is a CM-type, and $p_{K}\left(\alpha, \tau_{u}\right)^{2}=p_{K}\left(\alpha, \tau_{u}-\tau_{u} \rho\right)=p_{K}(\alpha, \Phi) p_{K}\left(\alpha, \Phi^{\prime}\right)^{-1}$. Therefore, to prove our proposition, it is sufficient to show that $p_{K}\left(\tau_{u}, \Phi\right)$ can be represented by a real number for every CM-type $\Phi$ and every $u \in \mathbf{a}$. Now we can choose $w_{0}$ so that $\tau_{v}\left(w_{0}\right)$ is pure imaginary for every $v \in$ a. Take $f=R_{u} \in \mathcal{A}_{\sigma_{u}}\left(F^{u}\right)$ as in Proposition 9.11 with $n=1$ in Case SP so that $f(w)$ is finite and nonzero. By Proposition 11.18, $f(w)$ represents $p_{K}\left(\tau_{u}, \Phi\right)$. Since $f$ is $F^{u}$-rational and the components of $w$ are pure imaginary, we see that $f(w) \in \mathbf{R}$. This proves our proposition.
11.20. Remark. (A) In the elliptic modular case, if $f \in \mathcal{A}_{0}(D)$ with a subfield $D$ of $\overline{\mathbf{Q}}$, then $(2 \pi i)^{-1} d f / d z$ belongs to $\mathcal{A}_{2}(D)$. Therefore, by the above lemma the value $(d f / d z)(w)$ for a CM-point $w$ belongs to $\pi p_{K}(\tau, \tau)^{2} \overline{\mathbf{Q}}$, where $K=\mathbf{Q}(w)$ and $\tau$ is the identity injection of $K$ into $\mathbf{C}$. We can naturally ask the nature of the derivatives of an element of $\mathcal{A}_{0}(\overline{\mathbf{Q}})$ in the general case. As one can easily imagine, this problem in the higher-dimensional case is highly nontrivial, and especially so when $\Gamma \backslash \mathcal{H}$ is compact. For example, suppose $F \neq \mathbf{Q}$ and $\sum_{v \in \mathbf{a}} m_{v} n_{v}=1$ in Case UB. Then $\operatorname{dim}(\mathcal{H})=1$, and so, given a nonconstant element $g$ of $\mathcal{A}_{0}(\overline{\mathbf{Q}})$, we have a well-defined derivation $d / d g$ of $\mathcal{A}_{0}(\overline{\mathbf{Q}})$. If $z$ is the variable on $\mathcal{H}$, then for every $f \in \mathcal{A}_{0}(\overline{\mathbf{Q}})$ we have $(d f / d z) /(d g / d z)=d f / d g \in \mathcal{A}_{0}(\overline{\mathbf{Q}})$. Therefore the value of $(d f / d z) /(d g / d z)$ at a CM-point $w$ is algebraic. But what exactly is the nature of the coset $(d f / d z)(w) \overline{\mathbf{Q}}$ ? Or more precisely, can we express it by means of the symbol $\mathfrak{p}(w)$ ? In fact, $(d f / d z)(w) \in \pi \mathfrak{p}_{v}(w) \overline{\mathbf{Q}}$, where $v$ is the element of a for which $m_{v} n_{v}=1$. In the next chapter we shall present a systematic treatment of differential operators acting on automorphic forms, especially such problems of arithmeticity in view. For example, the last fact on $(d f / d z)(w)$ is an easy special case of Proposition 14.5 below.
(B) There is a somewhat more intrinsic way of defining the symbol $\mathfrak{p}(w)$ than what was done in §11.4. In the setting of $\S 4.11$, we put $Y_{\mathbf{R}}=Y \otimes_{\mathbf{Q}} \mathbf{R}$ and $Y_{\mathbf{R}}^{u}=$ $\left\{x \in Y_{\mathbf{R}} \mid x x^{\rho}=1\right\}$, where $\rho$ is extended $\mathbf{R}$-linearly to $Y_{\mathbf{R}}$. Clearly $Y_{\mathbf{R}}^{u}$ can be identified with $\left\{\left(x_{\sigma}\right)_{\sigma \in J_{Y}} \in \mathbf{C}^{J_{Y}} \mid x_{\sigma} x_{\sigma \rho}=1\right\}$. Define an element $\mathfrak{q}_{w}$ of $Y_{\mathbf{R}}^{u}$ by

$$
\begin{equation*}
\mathfrak{q}_{w}=\left(p_{Y}(\sigma, \Phi)\right)_{\sigma \in J_{Y}} \tag{11.24}
\end{equation*}
$$

We extend $\psi_{v}$ and $\varphi_{v}$ of (4.37) to $\mathbf{R}$-linear maps of $Y_{\mathbf{R}}$ into $\mathbf{C}_{m_{v}}^{m_{v}}$ and $\mathbf{C}_{n_{v}}^{n_{v}}$; we then define $\lambda_{v}(h(\alpha), w)$ and $\mu_{v}(h(\alpha), w)$ for $\alpha \in Y_{\mathbf{R}}^{u}$ by (4.37), so that $M(h(\alpha), w)=$ $\left(\lambda_{v}(h(\alpha), w), \mu_{v}(h(\alpha), w)\right)_{v \in \mathbf{a}}$ is meaningful for $\alpha \in Y_{\mathbf{R}}^{u}$. Then (11.3a, b) and $(11.4 \mathrm{a}, \mathrm{b})$ show that $\lambda_{v}\left(h\left(\mathfrak{q}_{w}\right), w\right)=\psi_{v}\left(\mathfrak{q}_{w}\right)=\mathfrak{p}_{v \rho}(w)$ and $\mu_{v}\left(h\left(\mathfrak{q}_{w}\right), w\right)=$ $\varphi_{v}\left(\mathfrak{q}_{w}\right)=\mathfrak{p}_{v}(w)$. Therefore we can define $\mathfrak{p}(w)$ by

$$
\begin{equation*}
\mathfrak{p}(w)=M\left(h\left(\mathfrak{q}_{w}\right), w\right) \tag{11.25}
\end{equation*}
$$

(C) All arithmetic quotients treated in this book are associated with PEL-types. We can also construct canonical models of certain types of arithmetic quotients not associated with PEL-types, as we already noted it after the proof of Theorem 9.6, and then can define arithmeticity of automorphic forms in such cases. In each case we can speak of a CM-point, say $w$, obtained from a CM-algebra $Y$; then we have an element $\Phi$ of $I_{Y}$ and define the symbol $\mathfrak{p}(w)$ in terms of $p_{Y}(\alpha, \Phi)$ with $\alpha \in J_{Y}$ by natural analogues of (11.24) and (11.25). These are given in [S80, §4 resp. §5] when the group is a quaternion unitary group resp. an orthogonal group. One noteworthy aspect is that $\Phi$ is not necessarily a collection of CM-types for such groups.

## CHAPTER III

## ARITHMETIC OF DIFFERENTIAL OPERATORS and Nearly holomorphic functions

## 12. Differential operators on symmetric spaces

12.1. In this section we deal with two types of irreducible hermitian symmetric spaces $H$ of noncompact type, called Types A and C, which we already discussed in Section 3. To treat the space given in various different forms uniformly, we use symbols different from those of Section 3. For each type of $H$ we have a simple Lie group $G$, its maximal compact subgroup $K$ such that $G / K$ can be identified with $H$, and a complex vector space $T$ which can be identified with the holomorphic (or nonholomorphic) tangent space of $H$, on which the complexification $K^{c}$ of $K$ acts; $G$ is unitary for Type A and symplectic for Type C. The space $H$ can be presented as a "ball," or a "tube" under a certain condition. We refer to Types AB and CB if $H$ is a ball, and to Types AT and CT if it is a tube. If it is unnecessary to specify the distinction between a ball and a tube, we simply speak of Types A and C. In fact, we can treat all types of classical hermitian symmetric spaces of noncompact type by the same methods, but for simplicity we restrict our discussion to those two types; we refer the reader to [S94b] for a detailed treatment in the general case.

Now $G, K^{c}, H$, and $T$ for each type can be given explicitly in the table below. In each case we view $H$ as a subset of $T$; we also give positive definite hermitian matrices $\xi(z)$ and $\eta(z)$ defined for $z \in H$, which are closely connected with the Kähler structure of $H$ and the canonical factor of automorphy for the elements of $G$. We do not need $K$, nor any conceptual or geometric meaning of these objects, for the moment. Strictly speaking, $K^{c}$ is merely isomorphic, and not exactly equal, to the complexification of $K$. Here, however, we define $K^{c}$ formally as follows. Its connection with the maximal compact subgroup will be discussed in §A8.2.

$$
\begin{aligned}
& \text { Type AB: } \quad G=S U(m, n)=\left\{\alpha \in S L_{m+n}(\mathbf{C}) \mid \alpha^{*} I_{m, n} \alpha=I_{m, n}\right\}, \\
& I_{m, n}=\operatorname{diag}\left[1_{m},-1_{n}\right], \quad K^{c}=\left\{(a, b) \in G L_{m}(\mathbf{C}) \times G L_{n}(\mathbf{C}) \mid \operatorname{det}(a)=\operatorname{det}(b)\right\}, \\
& T=\mathbf{C}_{n}^{m}, \quad H=\left\{z \in T \mid z z^{*}<1_{m}\right\}, \\
& \xi(z)=1_{m}-\bar{z} \cdot{ }^{t} z, \quad \eta(z)=1_{n}-z^{*} z . \\
& \text { Type CB: } \quad G=S p(n, \mathbf{C}) \cap S U(n, n), \quad K^{c}=G L_{n}(\mathbf{C}) \text {, } \\
& T=\left\{\left.z \in \mathbf{C}_{n}^{n}\right|^{t} z=z\right\}, \quad H=\left\{z \in T \mid z^{*} z<1_{n}\right\}, \\
& \xi(z)=\eta(z)=1_{n}-z^{*} z .
\end{aligned}
$$

The space $H$ of Type CB is equivalent to a tube; however, $H$ of Type AB is equivalent to a tube if and only if $m=n$. We present here the explicit forms of the
objects associated to $H$ of tube form; we also present a real vector subspace $U$ of $T$ such that $T=U \otimes_{\mathbf{R}} \mathbf{C}$ and that $H=U+i P$ with a domain of positivity $P$ in $U$.

Type AT: $\quad G=S U\left(\eta_{n}\right)=\left\{\alpha \in S L_{2 n}(\mathbf{C}) \mid \alpha^{*} \eta_{n} \alpha=\eta_{n}\right\}, \quad \eta_{n}=\left[\begin{array}{cc}0 & -1_{n} \\ 1_{n} & 0\end{array}\right]$,

$$
\begin{aligned}
& K^{c}=\left\{(a, b) \in G L_{n}(\mathbf{C}) \times G L_{n}(\mathbf{C}) \mid \operatorname{det}(a)=\operatorname{det}(b)\right\} \\
& T=\mathbf{C}_{n}^{n}, \quad U=\left\{x \in \mathbf{C}_{n}^{n} \mid x^{*}=x\right\}, \quad H=\left\{z \in T \mid i\left(z^{*}-z\right)>0\right\} \\
&{ }^{t} \xi(z)=\eta(z)=i\left(z^{*}-z\right)
\end{aligned}
$$

Type CT: $\quad G=S p(n, \mathbf{R}), \quad K^{c}=G L_{n}(\mathbf{C})$,

$$
\begin{gathered}
T=\left\{\left.z \in \mathbf{C}_{n}^{n}\right|^{t} z=z\right\}, \quad U=\left\{\left.x \in \mathbf{R}_{n}^{n}\right|^{t} x=x\right\} \\
\xi(z)=\eta(z)=i(\bar{z}-z) .
\end{gathered}
$$

The space $H$ of Type CT, AT, or AB is exactly $\mathfrak{H}_{n}, \mathcal{H}_{n}$, or $\mathfrak{B}_{m, n}$ in Case SP, UT, or UB of $\S 3.1 ; H$ of Type CB is $\mathfrak{B}_{n}$ of (3.39). The functions $\xi$ and $\eta$ are the same as those of (3.9) and (3.10). The space $H$ of Type CB or CT is contained in $H$ of Type AB (with $m=n$ ) or AT; $\xi$ and $\eta$ for Type C are just restrictions of the corresponding functions for Type A.

In Section 3 we defined the action of $G$ on $H$ and also two factors of automorphy $\lambda_{\alpha}(z)$ and $\mu_{\alpha}(z)$; see (3.15), (3.16), and (3.17); $\lambda_{\alpha}(z)=\mu_{\alpha}(z)$ for Type C. We recall here a few basic formulas:

$$
\begin{gather*}
\lambda_{\alpha}(z)^{*} \xi(\alpha z) \lambda_{\alpha}(z)=\xi(z), \quad \mu_{\alpha}(z)^{*} \eta(\alpha z) \mu_{\alpha}(z)=\eta(z),  \tag{12.1a}\\
d(\alpha z)={ }^{t} \lambda_{\alpha}(z)^{-1} \cdot d z \cdot \mu_{\alpha}(z)^{-1} \quad(\alpha \in G) . \tag{12.1b}
\end{gather*}
$$

These were given in (3.19) and Lemma 3.4 (1). We also need a scalar factor of automorphy $j_{\alpha}(z)$ and a scalar-valued function $\delta$ defined in (3.20) and (3.21). Since our group $G$ of Type AB or AT is $S U(m, n)$ or $S U\left(\eta_{n}\right)$, the factors $\operatorname{det}(\alpha)$ and $\nu(\alpha)$ in the formulas in Section 3 are all equal to 1 . In particular we note that

$$
\begin{equation*}
\operatorname{det}\left(\lambda_{\alpha}(z)\right)=\operatorname{det}\left(\mu_{\alpha}(z)\right)=j_{\alpha}(z) \tag{12.2}
\end{equation*}
$$

We have $K^{c}=G L_{n}(\mathbf{C})$ for Type C. We view it as a subgroup of $G L_{n}(\mathbf{C}) \times G L_{n}(\mathbf{C})$ by the embedding $a \mapsto(a, a)$. We make a convention that $m$ means $n$ for Types $\mathrm{AT}, \mathrm{CT}$, and CB. Then, in all cases $K^{c}$ is a subgroup of $G L_{m}(\mathbf{C}) \times G L_{n}(\mathbf{C})$. Hereafter we shall write an element of $K^{c}$ in the form $(a, b) \in G L_{m}(\mathbf{C}) \times G L_{n}(\mathbf{C})$ with the understanding that $a=b$ for Type C. For example, $(\xi(z), \eta(z)) \in K^{c}$ for all types. (Notice that $\operatorname{det} \xi(z)=\operatorname{det} \eta(z)$ in all cases; see (3.13) for Type AB.) We let $K^{c}$ act on $T$ by $(a, b) u=a u \cdot{ }^{t} b$ for $(a, b) \in K^{c}$ and $u \in T$. For $T$ of Type C, this is the action of $G L_{n}(\mathbf{C})$ on the space of complex symmetric matrices defined by $u \mapsto a u \cdot{ }^{t} a$ for $a \in G L_{n}(\mathbf{C})$.
12.2. Throughout the rest of this section we denote by $G$ an algebraic group of the following types:

$$
\begin{array}{ll}
G=S U(\mathcal{T}) & \text { (Type AB), } \\
G=S U\left(\eta_{n}\right) & \text { (Type AT) }, \\
G=S p(n, F) & \text { (Type CT), } \\
G=S U\left(\eta_{n}\right) \cap S p(n, K) & \text { (Type CB). }
\end{array}
$$

Here $F$ is a totally real algebraic number field and $K$ is a totally imaginary quadratic extension of $F ; \mathcal{T}$ is an antihermitian element of $G L_{r}(K)$. Thus our group of Types $\mathrm{AB}, \mathrm{AT}$, and CT is exactly the group $G_{1}$ of $\S 8.2$, which is a subgroup of the group $G$ defined in $\S 3.5$ in Cases UB, UT, and SP. In this sense our notation is not consistent, but since in this section we exclusively consider the elements of $G_{1}$, we simply denote it by $G$; we do not consider larger groups like $G p(n, F)$ and $G U\left(\eta_{n}\right)$ in this and the next sections. Now we have

$$
\begin{equation*}
G_{\mathbf{a}}=\prod_{v \in \mathbf{a}} G_{v} \tag{12.3}
\end{equation*}
$$

and we see that each $G_{v}$ is either compact or a group belonging to the four types of $\S 1.1 ; G_{v}$ is not compact for $G$ of Types AT, CT and CB.

For each $v \in \mathbf{a}$ such that $G_{v}$ is not compact, we take the objects $H, T, K^{c}$, etc. associated with $G_{v}$, and denote them by $H_{v}, T_{v}, K_{v}^{c}$, etc. If $G_{v}$ is compact, that is, if $G_{v}=S U(m, n)$ with $m n=0$, then we put $K_{v}^{c}=S L_{m+n}(\mathbf{C}), T_{v}=\{0\}$, and let $H_{v}$ denote the space consisting of a single element 0 (see $\S 3.3$ ). We then put

$$
\begin{gather*}
\mathcal{H}=\prod_{v \in \mathbf{a}} H_{v}, \quad \mathfrak{K}_{0}=\prod_{v \in \mathbf{a}} K_{v}^{c},  \tag{12.4a}\\
\alpha z=\alpha(z)=\left(\alpha_{v} z_{v}\right)_{v \in \mathbf{a}}, \quad \Xi(z)=\left(\xi_{v}\left(z_{v}\right), \eta_{v}\left(z_{v}\right)\right)_{v \in \mathbf{a}}\left(\in \mathfrak{K}_{0}\right),  \tag{12.4b}\\
M_{\alpha}(z)=\left(\lambda\left(\alpha_{v}, z_{v}\right), \mu\left(\alpha_{v}, z_{v}\right)\right)_{v \in \mathbf{a}}\left(\in \mathfrak{K}_{0}\right) \tag{12.4c}
\end{gather*}
$$

for $z=\left(z_{v}\right)_{v \in \mathbf{a}} \in \mathcal{H}$ and $\alpha=\left(\alpha_{v}\right)_{v \in \mathbf{a}} \in G_{\mathbf{a}}$, where $\xi_{v}$ and $\eta_{v}$ denote the functions $\xi$ and $\eta$ defined on $H_{v}$. If $G_{v}$ is compact, we understand that the $v$-component of $\Xi$ is 1 , and the $v$-component of $M_{\alpha}$ is $\bar{\alpha}$ or $\alpha$ according to the convention of (32.4a, b). We also put

$$
\begin{equation*}
\mathbf{a}^{\prime}=\left\{v \in \mathbf{a} \mid G_{v} \text { is not compact }\right\} . \tag{12.5}
\end{equation*}
$$

In this section, by a representation $\{\sigma, W\}$ of a topological group $\mathcal{G}$, we mean a pair formed by a complex vector space $W$ of finite dimension and a continuous homomorphism $\sigma$ of $\mathcal{G}$ into $G L(W)$. We now take a representation $\{\rho, X\}$ of $\mathfrak{K}_{0}$ such that $\rho$ is complex analytic. Given $f \in C^{\infty}(\mathcal{H}, X)$ and $\alpha \in G_{\mathbf{a}}$, we define $f \|_{\rho} \alpha \in C^{\infty}(\mathcal{H}, X)$ by

$$
\begin{equation*}
\left(f \|_{\rho} \alpha\right)(z)=\rho\left(M_{\alpha}(z)\right)^{-1} f(\alpha z) \quad(z \in \mathcal{H}) \tag{12.6}
\end{equation*}
$$

In the previous sections we considered automorphic forms with respect to a representation $\{\omega, V\}$ of $\prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C})$. For Type C we have $\mathfrak{K}_{0}=\prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C})$, but for Type A, $\mathfrak{K}_{0}$ is smaller than $\prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C})$. Also, in this section we are taking the group $S U$ instead of $G U$ or $U$. For these reasons we use the letter $\rho$ instead of $\omega$. As we said in $\S 12.1$, we can develop the whole theory for an arbitray hermitian symmetric space, and our exposition, being almost axiomatic, can easily be generalized, as shown in [S94b]. In order to deal with such a general case, it is easier to consider semisimple groups, which is the reason why we consider $S U$ here. However, we shall later reinstate $\omega$, and take $\rho$ to be the restriction of $\omega$ to $\mathfrak{K}_{0}$. Anyway (12.6) is consistent with (5.6a).
12.3. Given a positive integer $p$ and finite-dimensional complex vector spaces $X$ and $Y$, we denote by $M l_{p}(Y, X)$ the vector space of all $\mathbf{C}$-multilinear maps of $Y \times \cdots \times Y$ ( $p$ copies) into $X$, and by $S_{p}(Y, X)$ the vector space of all homogeneous polynomial maps of $Y$ into $X$ of degree $p$. Thus $S_{1}(Y, X)=M l_{1}(Y, X)$, and it is the vector space of all C-linear maps of $Y$ into $X$. We put $S_{0}(Y, X)=M l_{0}(Y, X)=$
$X$ and $S_{p}(Y)=S_{p}(Y, \mathbf{C})$. We call an element $g$ of $M l_{p}(Y, X)$ symmetric if $g\left(y_{\pi(1)}, \ldots, y_{\pi(p)}\right)=g\left(y_{1}, \ldots, y_{p}\right)$ for every permutation $\pi$ of $\{1, \ldots, p\}$.

We are going to define differential operators on $\mathcal{H}$ with respect to each component $z_{v}$ of the variable point $z=\left(z_{v}\right)_{v \in \mathbf{a}}$ on $\mathcal{H}$ such that $G_{v}$ is not compact. Given such a $v, 0 \leq p \in \mathbf{Z}$, and a representation $\{\rho, X\}$ of $\mathfrak{K}_{0}$ as above, we define representations $\left\{\rho \otimes \tau_{v}^{p}, M l_{p}\left(T_{v}, X\right)\right\}$ and $\left\{\rho \otimes \sigma_{v}^{p}, M l_{p}\left(T_{v}, X\right)\right\}$ of $\mathcal{K}_{0}$ by

$$
\begin{align*}
& \quad\left[\left(\rho \otimes \tau_{v}^{p}\right)(a, b) h\right]\left(u_{1}, \ldots, u_{p}\right)=\rho(a, b) h\left({ }^{t} a_{v} u_{1} b_{v}, \ldots,{ }^{t} a_{v} u_{p} b_{v}\right)  \tag{12.7a}\\
& {\left[\left(\rho \otimes \sigma_{v}^{p}\right)(a, b) h\right]\left(u_{1}, \ldots, u_{p}\right)=\rho(a, b) h\left(a_{v}^{-1} u_{1}{ }^{t} b_{v}^{-1}, \ldots, a_{v}^{-1} u_{p}{ }^{t} b_{v}^{-1}\right)} \tag{12.7b}
\end{align*}
$$

for $(a, b) \in \mathfrak{K}_{0}, h \in M l_{p}\left(T_{v}, X\right)$, and $u_{i} \in T_{v}$. We use the same symbols $\rho \otimes \tau_{v}^{p}$ and $\rho \otimes \sigma_{v}^{p}$ for their restrictions to $S_{p}\left(T_{v}, X\right)$, and write them simply $\tau_{v}^{p}$ and $\sigma_{v}^{p}$ if $X=\mathbf{C}$ and $\rho$ is trivial. Here we use $(a, b)$ to denote an element of $\mathfrak{K}_{0}$ with $a=\left(a_{v}\right)_{v \in \mathbf{a}} \in \prod_{v \in \mathbf{a}} G L_{m_{v}}(\mathbf{C})$ and $b=\left(b_{v}\right)_{v \in \mathbf{a}} \in \prod_{v \in \mathbf{a}} G L_{n_{v}}(\mathbf{C})$, ignoring $a_{v}$ or $b_{v}$ according as $m_{v}=0$ or $n_{v}=0$.
12.4. Lemma. (1) $S_{p}(Y)$ is spanned by $g^{p}$ for all $g \in S_{1}(Y)$.
(2) Given $h \in S_{p}(Y, X)$, there is a unique symmetric element of $M l_{p}(Y, X)$, which we write $h_{*}$, such that $h(y)=h_{*}(y, \ldots, y)$.

Proof. We prove (1) by induction on $\operatorname{dim}(Y)$. Let $S_{p}^{\prime}(Y)$ be the subspace of $S_{p}(Y)$ spanned by $g^{p}$ for all $g \in S_{1}(Y)$. Clearly $S_{p}^{\prime}(Y)=S_{p}(Y)$ if $\operatorname{dim}(Y)=1$. Now we consider $\mathbf{C} \oplus Y$, and take a variable $(t, y) \in \mathbf{C} \oplus Y$. For $c \in \mathbf{C}$ and $f \in S_{1}(Y)$ we have $(c t+f(y))^{p}=\sum_{m=0}^{p}\binom{p}{m} c^{m} t^{m} f(y)^{p-m}$. From this we see that the functions of the form $t^{m} f(y)^{p-m}$ belong to $S_{p}^{\prime}(\mathbf{C} \oplus Y)$. By induction, $S_{p-m}^{\prime}(Y)=S_{p-m}(Y)$. Since $S_{p}(\mathbf{C} \oplus Y)$ is spanned by the functions $t^{m} g(y)$ with $g \in S_{p-m}(Y)$, we obtain (1) for $\mathbf{C} \oplus Y$. Clearly it is sufficient to prove (2) when $X=\mathbf{C}$. To prove the existence, by (1), we may assume that $h(y)=g(y)^{p}$ with $g \in S_{1}(Y)$. Then put $h_{*}\left(x_{1}, \ldots, x_{p}\right)=g\left(x_{1}\right) \cdots g\left(x_{p}\right)$. Clearly $h_{*}$ is symmetric and $h_{*}(y, \ldots, y)=h(y)$. We prove the uniqueness of $h_{*}$ by induction on $p$. Suppose $p>1$ and $f$ is a symmetric element of $M l_{p}(Y, \mathbf{C})$ such that $f(x, \ldots, x)=0$. Then

$$
0=f(x+y, \ldots, x+y)=\sum_{m=0}^{p}\binom{p}{m} f(\overbrace{x, \ldots, x}^{m}, \overbrace{y, \ldots, y}^{p-m}) .
$$

Viewing this as a polynomial function of $y$, we find that each term on the righthand side must vanish; in particular, $f(x, y, \ldots, y)=0$. Fixing $x$ and applying our induction to $f(x, y, \ldots, y)$, we obtain $f\left(x, x_{1}, \ldots, x_{p-1}\right)=0$ for arbitrary $x_{i}$, which completes the proof, since the case $p \leq 1$ is obvious.
12.5. We view $T_{v}$ as its own dual over $\mathbf{C}$ by the pairing $(u, v) \mapsto \operatorname{tr}\left({ }^{t} u v\right)$. Then, for $g \in S_{p}\left(T_{v}\right)$ and $h \in S_{p}\left(T_{v}, X\right)$ we put

$$
\begin{equation*}
[g, h]=\sum g_{*}\left(a_{\nu_{1}}, \ldots, a_{\nu_{p}}\right) h_{*}\left(b_{\nu_{1}}, \ldots, b_{\nu_{p}}\right) \tag{12.8}
\end{equation*}
$$

where $\left\{a_{\nu}\right\}_{\nu \in N}$ and $\left\{b_{\nu}\right\}_{\nu \in N}$ are dual bases of $T_{v}$, and ( $\nu_{1}, \ldots, \nu_{p}$ ) runs over $N^{p}$. Then $[g, h]$ is an element of $X$ determined independently of the choice of dual bases. From (12.7a, b) we easily obtain

$$
\begin{equation*}
\left[\sigma_{v}^{p}(\alpha) g,\left(\rho \otimes \tau_{v}^{p}\right)(\alpha) h\right]=\left[\tau_{v}^{p}(\alpha) g,\left(\rho \otimes \sigma_{v}^{p}\right)(\alpha) h\right]=\rho(\alpha)[g, h] \quad\left(\alpha \in \mathfrak{K}_{0}\right) \tag{12.9}
\end{equation*}
$$

In particular, taking $X=\mathbf{C}$, we can view $S_{p}\left(T_{v}\right)$ as its own dual by the pairing $(g, h) \mapsto[g, h]$, which is indeed nondegenerate because of a simple formula

$$
\begin{equation*}
[g, h]=g(x) \quad \text { if } \quad h(u)=\operatorname{tr}\left({ }^{t} x u\right)^{p} \text { with a fixed } x \in T_{v} . \tag{12.10}
\end{equation*}
$$

To prove this, we observe that $h_{*}\left(u_{1}, \ldots, u_{p}\right)=\prod_{i=1}^{p} \operatorname{tr}\left({ }^{t} x u_{i}\right)$; therefore $[g, h]=$ $\sum g_{*}\left(a_{\nu}, a_{\mu}, \cdots\right) \operatorname{tr}\left({ }^{t} x b_{\nu}\right) \operatorname{tr}\left({ }^{t} x b_{\mu}\right) \cdots$. Since $x=\sum_{\nu} \operatorname{tr}\left({ }^{t} x b_{\nu}\right) a_{\nu}$, we obtain (12.10).
12.6. Let us now recall some elementary facts on the polynomial representations of $G L_{n}(\mathbf{C})$. Every polynomial representation of $G L_{n}(\mathbf{C})$ can be decomposed into a direct sum of irreducible representations. The equivalence classes of irreducible polynomial representations of $G L_{n}(\mathbf{C})$ are in one-to-one correspondence with the ordered sets of integers $\left\{r_{1}, \ldots, r_{n}\right\}$ such that $r_{1} \geq \cdots \geq r_{n} \geq 0$. If $\{\sigma, W\}$ is such an irreducible representation corresponding to $\left\{r_{1}, \ldots, r_{n}\right\}$, then

$$
\operatorname{tr}\left\{\sigma\left(\operatorname{diag}\left[a_{1}, \ldots, a_{n}\right]\right)\right\}=a_{1}^{r_{1}} \cdots a_{n}^{r_{n}}+\sum_{s<r} \lambda_{s} a_{1}^{s_{1}} \cdots a_{n}^{s_{n}}
$$

for $a_{i} \in \mathbf{C}^{\times}$with $\lambda_{s} \in \mathbf{Z}$, where $<$ in the last sum is the lexicographic ordering. We call then $\sigma$ an irreducible representation of highest weight $\left\{r_{1}, \ldots, r_{n}\right\}$. A nonzero vector $q$ of $W$ is called an eigenvector of highest weight if $\sigma\left(\operatorname{diag}\left[a_{1}, \ldots, a_{n}\right]\right) q=$ $a_{1}^{r_{1}} \cdots a_{n}^{r_{n}} q$.

For $x \in \mathbf{C}_{n}^{m}$ and $1 \leq i \leq \operatorname{Min}(m, n)$ we denote by $\operatorname{det}_{i}(x)$ the determinant of the uper left $i^{2}$ entries of $x$. Let $R_{n}$ denote the subgroup of $G L_{n}(\mathbf{C})$ consisting of all the upper triangular matrices. Given an irreducible representation $\{\sigma, W\}$ of $G L_{n}(\mathbf{C})$, there is a unique common eigenvector $p$ of $\sigma\left(R_{n}\right)$ in $W$. If $\left\{r_{1}, \ldots, r_{n}\right\}$ is the highest weight of $\sigma$, then

$$
\begin{equation*}
\sigma(a) p=\prod_{i=1}^{n} \operatorname{det}_{i}(a)^{e_{i}} p \quad \text { for every } a \in R_{n} \tag{12.11}
\end{equation*}
$$

where $e_{i}=r_{i}-r_{i+1}$ for $i<n$ and $e_{n}=r_{n}$.
We are interested in the representations $\left\{S_{r}(T), \tau^{r}\right\}$ of a group $K^{c}$ defined in $\S 12.3$. We drop the subscript $v$. Thus we are looking at the objects as follows:
(Type A) $T=\mathbf{C}_{n}^{m},\left[\tau^{r}(a, b) h\right](u)=h\left({ }^{t} a u b\right)$ for $(a, b) \in \mathcal{K}=G L_{m}(\mathbf{C}) \times G L_{n}(\mathbf{C})$ and $h \in S_{r}(T)$;
(Type C) $T=\left\{x \in \mathbf{C}_{n}^{n},\left.\right|^{t} x=x\right\},\left[\tau^{r}(a) h\right](u)=h\left({ }^{t} a u a\right)$ for $a \in \mathcal{K}=G L_{n}(\mathbf{C})$ and $h \in S_{r}(T)$.
We have $\mathcal{K}=K^{c}$ for Type C , but $\mathcal{K}$ is larger than $K^{c}$ for Type $A$. However, clearly we can extend the representation $\tau^{r}$ to $\mathcal{K}$ as above, and the $\mathcal{K}$-irreducibility of a subspace of $S_{r}(T)$ is the same as that of $K^{c}$. Therefore we consider $\mathcal{K}$ instead of $K^{c}$. Also, for Type A we recall a well-known fact that every irreducible representation of $\mathcal{K}$ is of the form $\rho \otimes \sigma$ with irreducible representations $\rho$ of $G L_{m}(\mathbf{C})$ and $\sigma$ of $G L_{n}(\mathbf{C})$.
12.7. Theorem. Type A: Let $\rho$ and $\sigma$ be irreducible representations of $G L_{m}(\mathbf{C})$ and $G L_{n}(\mathbf{C})$, respectively. Then $\rho \otimes \sigma$ occurs in $\tau^{r}$ if and only if $\rho$ and $\sigma$ have the highest weights
$\left\{r_{1}, \ldots, r_{n}, 0, \ldots, 0\right\}$ and $\left\{r_{1}, \ldots, r_{n}\right\}$ when $m \geq n$,
$\left\{r_{1}, \ldots, r_{m}\right\}$ and $\left\{r_{1}, \ldots, r_{m}, 0, \ldots, 0\right\}$ when $n \geq m$,
with the $r_{i}$ such that $r_{1}+\cdots+r_{\mu}=r$ and $r_{\mu} \geq 0$, where $\mu=\operatorname{Min}(m, n)$. Such a $\rho \otimes \sigma$ has multiplicity one in $\tau^{r}$, and the corresponding irreducible subspace of $S_{r}(T)$ contains a polynomial $p$ defined by

$$
p(x)=\prod_{i=1}^{\mu} \operatorname{det}_{i}(x)^{e_{i}} \quad\left(x \in T=\mathbf{C}_{n}^{m}\right)
$$

as an eigenvector of highest weight with respect to both $\rho$ and $\sigma$, where $e_{i}=$ $r_{i}-r_{i+1}$ for $i<\mu$ and $e_{\mu}=r_{\mu}$.

Type C: An irreducible representation $\sigma$ of $G L_{n}(\mathbf{C})$ of highest weight $\left\{r_{1}, \ldots\right.$, $\left.r_{n}\right\}$ occurs in $\tau^{r}$ if and only if all $r_{i}$ are even, $r_{n} \geq 0$, and $r_{1}+\cdots+r_{n}=2 r$. Such a $\sigma$ has multiplicity one in $\tau^{r}$, and the corresponding irreducible subspace of $S_{r}(T)$ contains a polynomial $p$ defined by

$$
p(x)=\prod_{i=1}^{n} \operatorname{det}_{i}(x)^{e_{i}} \quad(x \in T)
$$

as an eigenvector of highest weight, where $e_{i}=\left(r_{i}-r_{i+1}\right) / 2$ for $i<n$ and $e_{n}=r_{n} / 2$.

The decomposition of $S_{r}(T)$ into irreducible subspaces with highest weights as described above is due to L.-K. Hua; the highest weight vector was determined by K. D. Johnson. The proof, as well as references to these and other related investigations, can be found in [S84b].
12.8. Lemma. (1) Let $Z$ and $W$ be different irreducible subspaces of $S_{r}(T)$. Then $[f, g]=0$ for every $f \in Z$ and every $g \in W$.
(2) The form $(f, g) \mapsto[f, g]$ is nondegenerate on any $K^{c}$-stable subspace of $S_{r}(T)$.

Proof. Wc prove this for Type A; Type C can be handled in a similar and simpler way. Take $Z$ to be the space described in Theorem 12.7 and take $p \in$ $Z$ as in that theorem. Define similarly a function $q$ in $W$. By (12.9) we have $\left[\tau^{r}\left({ }^{t} a,{ }^{t} b\right) f, g\right]=\left[f, \tau^{r}(a, b) g\right]$ for $f, g \in S_{r}(T)$. Taking $a$ and $b$ to be diagonal, we see that $[p, q]=0$. Since $p$ and $q$ are eigenvectors of $R_{m} \times R_{n}$, from (12.9) we obtain $\left[\tau^{r}\left({ }^{t} R_{m},{ }^{t} R_{n}\right) p, q\right]=0$, and consequently $\left[\tau^{r}\left({ }^{t} R_{m} R_{m},{ }^{t} R_{n} R_{n}\right) p, q\right]=0$. Now ${ }^{t} R_{n} R_{n}=\left\{x \in G L_{n}(\mathbf{C}) \mid \operatorname{det}_{i}(x) \neq 0\right.$ for every $\left.i \leq n\right\}$, and hence ${ }^{t} R_{m} R_{m} \times$ ${ }^{t} R_{n} R_{n}$ is dense in $G L_{m}(\mathbf{C}) \times G L_{n}(\mathbf{C})$. Therefore we see that $[Z, q]=0$, which combined with (12.9) proves (1). Since the form $(f, g) \mapsto[f, g]$ is nondegenerate on $S_{r}(T)$ as we noted in $\S 12.5$, (2) follows immediately from (1).
12.9. Given $C^{\infty}$ manifolds $V$ and $W$, we denote by $C^{\infty}(V, W)$ the set of all $C^{\infty}$ maps of $V$ into $W$. For the most part we take $\mathcal{H}$ as $V$. In fact, we will have to consider the derivatives of $C^{\infty}$ functions which are defined only on an open dense subset of $\mathcal{H}$, such as meromorphic functions on $\mathcal{H}$. For simplicity, however, we state our definitions and formulas only for $C^{\infty}$ functions on $\mathcal{H}$, but we will apply them to functions of more general types, without any additional comment, as the validity of such applications is obvious in each case.

Since $T_{v}$ has a natural $\mathbf{R}$-structure, we can speak of an $\mathbf{R}$-rational basis of $T_{v}$ over C. Take any such basis $\left\{\varepsilon_{\nu}\right\}_{\nu \in N}$, and for $u \in T_{v}$ define $u_{\nu} \in \mathbf{C}$ by $u=\sum_{\nu \in N} u_{\nu} \varepsilon_{\nu}$. We also put $z_{v}=\sum_{\nu \in N} z_{v \nu} \varepsilon_{\nu}$ with $z_{v \nu} \in \mathbf{C}$ for the variable $z_{v}$ on $H_{v}$. Then, for $f \in C^{\infty}(\mathcal{H}, X)$ we define $D_{v} f, \bar{D}_{v} f, C_{v} f, E_{v} f \in C^{\infty}\left(\mathcal{H}, S_{1}\left(T_{v}, X\right)\right)$ by

$$
\begin{array}{ll}
\left(D_{v} f\right)(u)=\sum_{\nu \in N} u_{\nu} \partial f / \partial z_{v \nu}, \quad\left(\bar{D}_{v} f\right)(u)=\sum_{\nu \in N} u_{\nu} \partial f / \partial \bar{z}_{v \nu}, \\
\left(C_{v} f\right)(u)=\left(D_{v} f\right)\left({ }^{t} \xi_{v} u \eta_{v}\right), \quad\left(E_{v} f\right)(u)=\left(\bar{D}_{v} f\right)\left(\xi_{v} u \cdot{ }^{t} \eta_{v}\right) \tag{12.12b}
\end{array}
$$

for $u \in T_{v}$. These are independent of the choice of $\left\{\varepsilon_{\nu}\right\}_{\nu \in N}$. The last two formulas can be written $C_{v} f=\tau_{v}^{1}(\Xi) D_{v} f$ and $E_{v} f=\sigma_{v}^{1}\left(\Xi^{-1}\right) \bar{D}_{v} f$ with $\Xi$ of (12.4b), and the first two are equivalent to the expression

$$
\begin{equation*}
d f=\sum_{v \in \mathbf{a}}\left(D_{v} f\right)\left(d z_{v}\right)+\sum_{v \in \mathbf{a}}\left(\bar{D}_{v} f\right)\left(d \bar{z}_{v}\right) . \tag{12.13}
\end{equation*}
$$

Notice that $E_{v} f=0$ if and only if $f$ is holomorphic in $z_{v}$. Substituting $f \circ \alpha$ for $f$ in (12.13) and employing (12.1b), we find that

$$
\begin{align*}
& \left(\left(D_{v} f\right) \circ \alpha\right)(u)=D_{v}(f \circ \alpha)\left({ }^{t} \lambda\left(\alpha_{v}, z_{v}\right) u \mu\left(\alpha_{v}, z_{v}\right)\right),  \tag{12.14a}\\
& \left(\left(\bar{D}_{v} f\right) \circ \alpha\right)(u)=\bar{D}_{v}(f \circ \alpha)\left(\lambda\left(\alpha_{v}, z_{v}\right)^{*} u \overline{\mu\left(\alpha_{v} z_{v}\right)}\right) \tag{12.14b}
\end{align*}
$$

for $u \in T_{v}$ and $\alpha \in G_{\mathbf{a}}$.
We also define $D_{v}^{e} f, \bar{D}_{v}^{e} f, C_{v}^{e} f$, and $E_{v}^{e} f$ for $0 \leq e \in \mathbf{Z}$ by

$$
\begin{align*}
& D_{v}^{e} f=D_{v} D_{v}^{e-1} f, \quad \bar{D}_{v}^{e} f=\bar{D}_{v} \bar{D}_{v}^{e-1} f, \quad D_{v}^{0} f=\bar{D}_{v}^{0} f=f,  \tag{12.15a}\\
& C_{v}^{e} f=C_{v} C_{v}^{e-1} f, \quad E_{v}^{e} f=E_{v} E_{v}^{e-1} f, \quad C_{v}^{0} f=E_{v}^{0} f=f . \tag{12.15b}
\end{align*}
$$

These have values in $M l_{e}\left(T_{v}, X\right)$ in the sense that
(12.15c) $\quad\left(A^{e} f\right)\left(u_{1}, \ldots, u_{e}\right)=A\left\{\left(A^{e-1} f\right)\left(u_{1}, \ldots, u_{e-1}\right)\right\}\left(u_{e}\right) \quad\left(u_{i} \in T_{v}\right)$,
where $A$ is $D_{v}, \bar{D}_{v}, C_{v}$ or $E_{v}$. Clearly $D_{v}^{e} f$ and $\bar{D}_{v}^{e} f$ have symmetric elements of $M l_{e}\left(T_{v}, X\right)$ as their values; the same is true for $C_{v}^{e} f$ and $E_{v}^{e} f$ as will be shown in $\S 13.9$. Therefore we can view them as elements of $C^{\infty}\left(\mathcal{H}, S_{e}\left(T_{v}, X\right)\right)$. For example, we have

$$
\begin{equation*}
\left(D_{v}^{p} f\right)(u)=p!h(u) \quad \text { if } f(z)=h\left(z_{v}\right) \text { with } h \in S_{p}\left(T_{v}\right) . \tag{12.16}
\end{equation*}
$$

Indeed, $h\left(z_{v}\right)=h_{*}\left(z_{v}, \ldots, z_{v}\right)$ with $h_{*}$ as in Lemma 12.4 (2), and hence $\left(D_{v} f\right)(u)$ $=p h_{*}\left(u, z_{v}, \ldots, z_{v}\right)$, and similarly $\left(D_{v}^{2} f\right)(u, v)=p(p-1) h_{*}\left(u, v, z_{v}, \ldots, z_{v}\right)$. Repeating this procedure, we obtain (12.16).

We now define $D_{\rho, v}^{e} f \in C^{\infty}\left(\mathcal{H}, S_{e}\left(T_{v}, X\right)\right)$ by

$$
\begin{equation*}
D_{\rho, v}^{e} f=\left(\rho \otimes \tau_{v}^{e}\right)(\Xi)^{-1} C_{v}^{e}[\rho(\Xi) f] . \tag{12.17}
\end{equation*}
$$

In particular, since $C_{v} f=\tau_{v}^{1}(\Xi) D_{v} f$, writing $D_{\rho, v}=D_{\rho, v}^{1}$, we have

$$
\begin{equation*}
\left(D_{\rho, v} f\right)(u)=\rho(\Xi)^{-1} D_{v}[\rho(\Xi) f](u) \quad\left(u \in T_{v}\right) \tag{12.18}
\end{equation*}
$$

12.10. Proposition. (1) $D_{\rho, v}^{e+1}=D_{\rho \otimes \tau_{v}^{e}, v} D_{\rho, v}^{e}=D_{\rho \otimes \tau_{v}, v}^{e} D_{\rho, v}$.
(2) $\quad D_{\rho, v}^{e}\left(f \|_{\rho} \alpha\right)=\left(D_{\rho, v}^{e} f\right)\left\|_{\rho \otimes \tau_{v}^{e}} \alpha, \quad E_{v}^{e}\left(f \|_{\rho} \alpha\right)=\left(E_{v}^{e} f\right)\right\|_{\rho \otimes \sigma_{v}^{e}} \alpha . \quad\left(\alpha \in G_{\mathbf{a}}\right)$.

Proof. Identifying $\rho \otimes \tau_{v}^{e} \otimes \tau_{v}$ with $\rho \otimes \tau_{v}^{e+1}$, by (12.17) and (12.18) we have

$$
\begin{aligned}
D_{\rho \otimes \tau_{v}^{e} \cdot v}\left(D_{\rho, v}^{e} f\right) & =\left(\rho \otimes \tau_{v}^{e+1}\right)(\Xi)^{-1} C_{v}\left[\left(\rho \otimes \tau_{v}^{e}\right)(\Xi)\left(D_{\rho, v}^{e} f\right)\right] \\
& =\left(\rho \otimes \tau_{v}^{e+1}\right)(\Xi)^{-1} C_{v} C_{v}^{e}[\rho(\Xi) f]=D_{\rho, v}^{e+1} f,
\end{aligned}
$$

which proves the first part of (1). The second part can be proved similarly. Now for $\alpha \in G_{\mathbf{a}}$ and $u \in T_{v}$ we have, employing (12.1a) and (12.14a),

$$
\begin{aligned}
& D_{\rho . v}\left(f \|_{\rho} \alpha\right)(u)=\rho(\Xi)^{-1} D_{v}\left[\rho\left(\Xi M_{\alpha}^{-1}\right)(f \circ \alpha)\right](u) \\
& \quad=\rho(\Xi)^{-1} D_{v}\left[\rho\left(M_{\alpha}^{*}\right)(\{\rho(\Xi) f\} \circ \alpha)\right](u)=\rho\left(\Xi^{-1} M_{\alpha}^{*}\right) D_{v}[(\rho(\Xi) f) \circ \alpha](u) \\
& \quad=\rho\left(M_{\alpha}^{-1}(\Xi \circ \alpha)^{-1}\right)\left\{D_{v}(\rho(\Xi) f) \circ \alpha\right\}\left({ }^{t} \lambda\left(\alpha_{v}, z_{v}\right)^{-1} u \mu\left(\alpha_{v}, z_{v}\right)^{-1}\right) \\
& \quad=\rho\left(M_{\alpha}^{-1}\right)\left\{\left(D_{\rho . v} f\right) \circ \alpha\right\}\left({ }^{t} \lambda\left(\alpha_{v}, z_{v}\right)^{-1} u \mu\left(\alpha_{v}, z_{v}\right)^{-1}\right) .
\end{aligned}
$$

Observing that the last quantity can be written $\left[\left(D_{\rho . v} f\right) \|_{\rho \& \tau_{v}} \alpha\right](u)$, we obtain the first formula of (2) for $e=1$. Then the general case can be proved by induction on $e$, by virtue of (1). As for $E_{v}$, we have similarly, by (12.14b) and (12.1a),

$$
\left.\begin{array}{l}
E_{v}\left(f \|_{\rho} \alpha\right)(u)=E_{v}\left\{\rho\left(M_{\alpha}^{-1}\right)(f \circ \alpha)\right\}(u)=\rho\left(M_{\alpha}\right)^{-1} \bar{D}_{v}(f \circ \alpha)\left(\xi_{v} u \cdot{ }^{t} \eta_{v}\right) \\
\quad=\rho\left(M_{\alpha}\right)^{-1}\left(\left(\bar{D}_{v} f\right) \circ \alpha\right)\left({ }^{t} \overline{\lambda\left(\alpha_{v}, z_{v}\right.}\right)^{-1} \xi_{v} u \cdot{ }^{t} \eta_{v} \overline{\mu\left(\alpha_{v}, z_{v}\right)}
\end{array}\right), ~=\rho\left(M_{\alpha}\right)^{-1}\left(\left(\bar{D}_{v} f\right) \circ \alpha\right)\left(\left(\xi_{v} \circ \alpha\right) \lambda\left(\alpha_{v}, z_{v}\right) u \cdot{ }^{t} \mu\left(\alpha_{v}, z_{v}\right)\left({ }^{t} \eta_{v} \circ \alpha\right)\right) .
$$

This proves the second formula of (2) for $e=1$. The general case can be proved by induction on $e$.
12.11. The representation $\tau_{v}^{p}$ or $\sigma_{v}^{p}$ of $\mathfrak{K}_{0}$ on $S_{p}\left(T_{v}\right)$ is essentially that of $K_{v}^{c}$. By Theorem 12.7 it is the direct sum of irreducible representations, and each irreducible constituent has multiplicity one. (Since the $\tau_{v}^{p}$-irreducibility is the same as the $\sigma_{v}^{p}{ }^{-}$ irreducibility, we shall simply speak of an irreducible subspace of $S_{p}\left(T_{v}\right)$.) Thus, for each $\mathfrak{K}_{0}$-stable subspace $Z$ of $S_{p}\left(T_{v}\right)$, we can define the projection map $\varphi_{Z}$ of $S_{p}\left(T_{v}\right)$ onto $Z$. Now we can identify $S_{p}\left(T_{v}, X\right)$ with $S_{p}\left(T_{v}\right) \otimes X$ by the rule

$$
\begin{equation*}
(h \otimes x)(u)=h(u) x \quad \text { for } h \in S_{p}\left(T_{v}\right), x \in X, \text { and } u \in T_{v} \tag{12.19}
\end{equation*}
$$

(This justifies the notation $\rho \otimes \tau_{v}^{p}$ and $\rho \otimes \sigma_{v}^{p}$.) Using the same symbol $\varphi_{Z}$ for the map $h \otimes x \mapsto\left(\varphi_{Z} h\right) \otimes x$ of $S_{p}\left(T_{v}\right) \otimes X$ to $Z \otimes X$, we define $D_{\rho}^{Z} f, E^{Z} f \in$ $C^{\infty}(\mathcal{H}, Z \otimes X)$ by

$$
\begin{equation*}
D_{\rho}^{Z} f=\varphi_{Z} D_{\rho, v}^{p} f, \quad E^{Z} f=\varphi_{Z} E_{v}^{p} f \tag{12.20}
\end{equation*}
$$

Let $\tau_{Z}$ and $\sigma_{Z}$ denote the restrictions of $\tau_{v}^{p}$ and $\sigma_{v}^{p}$ to $Z$. Then $\rho \otimes \tau_{Z}$ and $\rho \otimes \sigma_{Z}$ are the restrictions of $\rho \otimes \tau_{v}^{p}$ and $\rho \otimes \sigma_{v}^{p}$ to $Z \otimes X$. Then $\left(\varphi_{Z} f\right) \|_{\rho \otimes \tau_{Z}} \alpha=$ $\varphi_{Z}\left(f \|_{\rho \otimes \tau_{v}^{e}} \alpha\right)$ and $\left(\varphi_{Z} f\right) \|_{\rho \otimes \sigma_{Z}} \alpha=\varphi_{Z}\left(f \|_{\rho \otimes \sigma_{v}^{e}} \alpha\right)$. Therefore from Proposition 12.10 (2) we obtain, for every $\alpha \in G_{\mathbf{a}}$,

$$
\begin{equation*}
D_{\rho}^{Z}\left(f \|_{\rho} \alpha\right)=\left(D_{\rho}^{Z} f\right)\left\|_{\rho \otimes \tau_{Z}} \alpha, \quad E^{Z}\left(f \|_{\rho} \alpha\right)=\left(E^{Z} f\right)\right\|_{\rho \otimes \sigma_{Z}} \alpha \tag{12.21}
\end{equation*}
$$

By Lemma $12.8, Z$ is its own dual with respect to the bilinear form of (12.8). Therefore we can identify $Z \otimes X$ with $S_{1}(Z, X)$ by the rule

$$
\begin{equation*}
(\omega \otimes x)(\zeta)=[\zeta, \omega] x \quad \text { for } \omega, \zeta \in Z \text { and } x \in X \tag{12.22}
\end{equation*}
$$

Then $\varphi_{Z}$ as a $\operatorname{map} S_{p}\left(T_{v}, X\right) \rightarrow S_{1}(Z, X)$ can be given by $\left(\varphi_{Z} g\right)(\zeta)=[\zeta, g]$ for $g \in S_{p}\left(T_{v}, X\right)$ and $\zeta \in Z$, and hence $D_{\rho}^{Z} f$ and $E^{Z} f$ as $S_{1}(Z, X)$-valued functions can be given by

$$
\begin{equation*}
\left(D_{\rho}^{Z} f\right)(\zeta)=\left[\zeta, D_{\rho, v}^{p} f\right], \quad\left(E^{Z} f\right)(\zeta)=\left[\zeta, E_{v}^{p} f\right] \quad(\zeta \in Z) \tag{12.23}
\end{equation*}
$$

The symbols $\rho \otimes \tau_{Z}$ and $\rho \otimes \sigma_{Z}$ as representations on the space $S_{1}(Z, X)$ can be given by

$$
\begin{align*}
& {\left[\left(\rho \otimes \tau_{Z}\right)(a, b) h\right](\zeta)=\rho(a, b) h\left(\tau_{Z}\left({ }^{t} a_{v},{ }^{t} b_{v}\right) \zeta\right)}  \tag{12.24a}\\
& {\left[\left(\rho \otimes \sigma_{Z}\right)(a, b) h\right](\zeta)=\rho(a, b) h\left(\sigma_{Z}\left({ }^{t} a_{v},{ }^{t} b_{v}\right) \zeta\right)} \tag{12.24b}
\end{align*}
$$

for $h \in S_{1}(Z, X), \zeta \in Z$, and $(a, b) \in \mathfrak{K}_{0}$ as in (12.7a, b). Take, for example, $Z=\mathbf{C} \psi \subset S_{e n}(T)$ with $\psi(u)=\operatorname{det}(u)^{e}, 0 \leq e \in \mathbf{Z}$, assuming that $m=n$ for Type AB. Then, from (12.21) and (12.24a, b) we obtain

$$
\begin{equation*}
D_{\rho}^{Z}\left(f \|_{\rho} \alpha\right)(\psi)=\left[\left(D_{\rho}^{Z} f\right)(\psi)\right]\left\|_{\rho^{\prime}} \alpha, \quad E^{Z}\left(f \|_{\rho} \alpha\right)(\psi)=\left[\left(E^{Z} f\right)(\psi)\right]\right\|_{\rho^{\prime \prime}} \alpha \tag{12.24c}
\end{equation*}
$$

$$
\text { where } \rho^{\prime}(a, b)=\operatorname{det}(a)^{e} \operatorname{det}(b)^{e} \rho(x, y) \text { and } \rho^{\prime \prime}(a, b)=\operatorname{det}(a)^{-e} \operatorname{det}(b)^{-e} \rho(x, y)
$$

12.12. Fixing $v \in \mathbf{a}$ and taking $\mathbf{R}$-rational dual bases $\left\{\varepsilon_{\nu}\right\}_{\nu \in N}$ and $\left\{\varepsilon_{\nu}^{\prime}\right\}_{\nu \in N}$ of $T_{v}$ over $\mathbf{C}$, put $z_{v}=\sum_{\nu \in N} z_{v \nu} \varepsilon_{\nu}$ as in $\S 12.9$, and

$$
\begin{equation*}
\mathcal{D}_{v}=\sum_{\nu \in N} \varepsilon_{\nu}^{\prime} \partial / \partial z_{v \nu}, \quad \overline{\mathcal{D}}_{v}=\sum_{\nu \in N} \varepsilon_{\nu}^{\prime} \partial / \partial \bar{z}_{v \nu} \tag{12.25}
\end{equation*}
$$

These are independent of the choice of bases. For $g \in S_{p}\left(T_{v}\right)$ we define $g\left(\mathcal{D}_{v}\right)$ by

$$
\begin{equation*}
g\left(\mathcal{D}_{v}\right)=g_{*}\left(\mathcal{D}_{v}, \ldots, \mathcal{D}_{v}\right)=\sum g_{*}\left(\varepsilon_{\nu_{1}}^{\prime}, \ldots, \varepsilon_{\nu_{p}}^{\prime}\right) \partial^{p} / \partial z_{v \nu_{1}} \cdots \partial z_{v \nu_{p}} \tag{12.26}
\end{equation*}
$$

where $\left(\nu_{1}, \ldots, \nu_{p}\right)$ runs over $N^{p}$, and define $g\left(\overline{\mathcal{D}}_{v}\right)$ similarly. Now we have

$$
\begin{align*}
& g\left(\mathcal{D}_{v}\right) f=\left[g, D_{v}^{p} f\right] \quad \text { for every } \quad f \in C^{\infty}(\mathcal{H}, X)  \tag{12.27}\\
& g\left(\mathcal{D}_{v}\right) h=p![g, h] \quad \text { for every } h \in S_{p}\left(T_{v}\right) \tag{12.28}
\end{align*}
$$

with $[g, h]$ of (12.8). In (12.27) we view $D_{v}^{p} f$ as $S_{p}\left(T_{v}, X\right)$-valued; in (12.28) we view $h$ as a function on $\mathcal{H}$ by putting $h(z)=h\left(z_{v}\right)$. To prove these, suppress the subscript $v$. Since $(D f)\left(\varepsilon_{\nu}\right)=\partial f / \partial z_{\nu}$, we have $\left(D^{p} f\right)\left(\varepsilon_{\nu_{1}}, \ldots, \varepsilon_{\nu_{p}}\right)=$ $\partial^{p} f / \partial z_{\nu_{1}} \cdots \partial z_{\nu_{p}}$, which combined with (12.8) and (12.26) proves (12.27). Then (12.28) follows from (12.16) and (12.27).

For example, since $\left[D_{v}^{p} \exp \left(\operatorname{tr}\left({ }^{t} x z_{v}\right)\right)\right](u)=\operatorname{tr}\left({ }^{t} x u\right)^{p} \exp \left(\operatorname{tr}\left({ }^{t} x z_{v}\right)\right)$ for $u, x \in$ $T_{v}$, from (12.10) and (12.27) we obtain

$$
\begin{equation*}
g\left(\mathcal{D}_{v}\right) \exp \left(\operatorname{tr}\left({ }^{t} x z_{v}\right)\right)=g(x) \exp \left(\operatorname{tr}\left({ }^{t} x z_{v}\right)\right) \tag{12.29}
\end{equation*}
$$

In particular, for an element $g$ of $S_{n k}\left(T_{v}\right)$ given by $g(u)=\operatorname{det}(u)^{k}$, we denote the operator $g\left(\mathcal{D}_{v}\right)$ by $\operatorname{det}\left(\mathcal{D}_{v}\right)^{k}$.

In the following theorem, we drop the subscript $v$ for simplicity.
12.13. Theorem. Let $Z$ be the irreducible subspace of $S_{r}(T)$ described in Theorem 12.7 and $r_{i}$ be the integers in that theorem; let $\zeta \in Z, s \in \mathbf{C}$, and

$$
L= \begin{cases}\left\{(c, d) \in \mathbf{C}_{m}^{n} \times \mathbf{C}_{n}^{n} \left\lvert\, \operatorname{rank}\left[\begin{array}{ll}
c & d
\end{array}\right]=n\right.\right\} \\
\left\{(c, d) \in \mathbf{C}_{n}^{n} \times \mathbf{C}_{n}^{n} \mid \operatorname{rank}[c\right. & \left.d]=n, c \cdot{ }^{t} d=d \cdot{ }^{t} c\right\}\end{cases}
$$

Then, for any fixed $(c, d) \in L$ and any fixed branch of $\operatorname{det}(c z+d)^{s}$ in an open subset of $T$ on which $\operatorname{det}(c z+d)^{s}$ is meaningful, we have

$$
\begin{aligned}
& \zeta(\mathcal{D}) \operatorname{det}(c z+d)^{s}=\psi_{Z}(s) \operatorname{det}(c z+d)^{s} \zeta\left({ }^{t} c \cdot{ }^{t}(c z+d)^{-1}\right) \\
& \text { with } \quad \psi_{Z}(s)= \begin{cases}\prod_{h=1}^{\mu} \prod_{i=1}^{r_{h}}(s-i+h), \quad \mu=\operatorname{Min}(m, n) \\
\prod_{h=1}^{n} \prod_{i=1}^{r_{h} / 2}\left(s-i+\frac{h+1}{2}\right) & \text { (Type A), }\end{cases}
\end{aligned}
$$

For the proof, see [S84b, Theorem 4.3].
12.14. The notation being as in $\S 12.11$, we define a contraction operator $\theta$ : $Z \otimes Z \otimes X \rightarrow X$ by $\theta(\zeta \otimes \omega \otimes x)=[\zeta, \omega] x$. This as a map $S_{1}\left(Z, S_{1}(Z, X)\right) \rightarrow X$ can be given by

$$
\begin{equation*}
\theta h=\sum_{\mu} h\left(\zeta_{\mu}, \omega_{\mu}\right) \quad \text { for } \quad h \in S_{1}\left(Z, S_{1}(Z, X)\right) \tag{12.30a}
\end{equation*}
$$

with bases $\left\{\zeta_{\mu}\right\}$ and $\left\{\omega_{\mu}\right\}$ of $Z$ such that $\left[\zeta_{\mu}, \omega_{\nu}\right]=\delta_{\mu \nu}$. From (12.9) we obtain
$(12.30 \mathrm{~b}) \quad \theta \circ\left(\rho \otimes \tau_{Z} \otimes \sigma_{Z}\right)(a, b)=\theta \circ\left(\rho \otimes \sigma_{Z} \otimes \tau_{Z}\right)(a, b)=\rho(a, b) \circ \theta$ for every $(a, b) \in \mathfrak{K}_{0}$.

If $g \in C^{\infty}\left(\mathcal{H}, S_{1}(Z, X)\right)$, then $D_{\rho \otimes \sigma_{Z}}^{Z} g$ and $E^{Z} g$ have values in $S_{1}\left(Z, S_{1}(Z, X)\right)$, so that $\theta D_{\rho \otimes \sigma_{Z}}^{Z} g$ and $\theta E^{Z} g$ are meaningful as $X$-valued functions. In particular, for $f \in C^{\infty}(\mathcal{H}, X)$ the symbols $\theta D_{\rho \otimes \sigma_{Z}}^{Z} E^{Z} f$ and $\theta E^{Z} D_{\rho}^{Z} f$ are elements of $C^{\infty}(\mathcal{H}, X)$. We then put

$$
\begin{equation*}
L_{\rho}^{Z} f=(-1)^{p} \theta D_{\rho \otimes \sigma_{Z}}^{Z} E^{Z} f, \quad M_{\rho}^{Z} f=(-1)^{p} \theta E^{Z} D_{\rho}^{Z} f \tag{12.31}
\end{equation*}
$$

Then, by (12.21) and (12.30b), for every $\alpha \in G_{\mathrm{a}}$ we have

$$
\begin{equation*}
L_{\rho}^{Z}\left(f \|_{\rho} \alpha\right)=\left(L_{\rho}^{Z} f\right)\left\|_{\rho} \alpha, \quad M_{\rho}^{Z}\left(f \|_{\rho} \alpha\right)=\left(M_{\rho}^{Z} f\right)\right\|_{\rho} \alpha \tag{12.32}
\end{equation*}
$$

If $Z=S_{1}\left(T_{v}\right)$ and this is identified with $T_{v}$, then the map $\theta$ can be viewed as a $\operatorname{map} S_{1}\left(T_{v}, S_{1}\left(T_{v}, X\right)\right) \rightarrow X$, and is given by

$$
\begin{equation*}
\theta h=\sum_{\nu \in N} h\left(a_{\nu}, b_{\nu}\right) \quad \text { for } \quad h \in S_{1}\left(T_{v}, S_{1}\left(T_{v}, X\right)\right) \tag{12.33}
\end{equation*}
$$

with dual bases $\left\{a_{\nu}\right\}$ and $\left\{b_{\nu}\right\}$ of $T_{v}$.
Given a congruence subgroup $\Gamma$ of $G$, we denote by $C_{\rho}(\Gamma)$ the set of all $f \in$ $C^{\infty}(\mathcal{H}, X)$ such that $f \|_{\rho} \gamma=f$ for every $\gamma \in \Gamma$, and by $C_{\rho}$ the union of $C_{\rho}(\Gamma)$ for all such $\Gamma$. Now $X$ has an inner product $\langle x, y\rangle_{X}$ which is $\mathbf{C}$-linear in $y$ and C-antilinear in $x$ and which satisfies

$$
\begin{equation*}
\langle x, \rho(a, b) y\rangle_{X}=\left\langle\rho\left(a^{*}, b^{*}\right) x, y\right\rangle_{X} \tag{12.34}
\end{equation*}
$$

for $(a, b) \in \mathfrak{K}_{0}$, where $\left(a^{*}, b^{*}\right)=\left(a_{v}^{*}, b_{v}^{*}\right)_{v \in \mathbf{a}}$. This will be proven after (12.35b). Then, for $f, g \in C_{\rho}$ we define their inner product $\langle f, g\rangle$ by

$$
\begin{equation*}
\langle f, g\rangle=\mu(\Phi)^{-1} \int_{\Phi}\langle f(z), \rho(\Xi(z)) g(z)\rangle_{X} d \mu(z) \tag{12.35a}
\end{equation*}
$$

whenever the integral is convergent, where $d \mu(z)$ is a fixed $G_{\mathbf{a}}$-invariant measure on $\mathcal{H}, \mu(\Phi)=\int_{\Phi} d \mu(z)$, and $\Phi=\Gamma \backslash \mathcal{H}$ with $\Gamma$ such that $f, g \in C_{\rho}(\Gamma)$. The inner product is independent of the choice of $\Gamma$, and

$$
\begin{equation*}
\left\langle f\left\|_{\rho} \alpha, g\right\|_{\rho} \alpha\right\rangle=\langle f, g\rangle \quad \text { for every } \quad \alpha \in G . \tag{12.35b}
\end{equation*}
$$

Let us now prove the existence of $\langle,\rangle_{X}$ satisfying (12.34). In view of the structure of $\mathfrak{K}_{0}$, it is sufficient to prove that if $(Y, \sigma)$ is an irrreducible polynomial representation of $G L_{n}(\mathbf{C})$, then $Y$ has an inner product $\langle,\rangle_{Y}$ such that $\langle x, \sigma(a) y\rangle_{Y}=\left\langle\sigma\left(a^{*}\right) x, y\right\rangle_{Y}$ for $a \in G L_{n}(\mathbf{C})$. Now $(Y, \sigma)$ is a direct summand of $(W, \omega)$, where $W=\mathbf{C}^{n} \otimes \cdots \otimes \mathbf{C}^{n}$ and $\omega(x)=x \otimes \cdots \otimes x$ with $\mathbf{C}^{n}$ and $x$ repeated $m$ times for some $m$. Clearly $W$ has an inner product with the required property with respect to $\omega$, and we only have to restrict it to $Y$.
12.15. Theorem. Let $Z$ be an irreducible subspace of $S_{p}\left(T_{v}\right)$. Then $Z \otimes X$ has an inner product satisfying (12.34) with $\left(Z \otimes X, \rho \otimes \tau_{Z}\right)$ and $\left(Z \otimes X, \rho \otimes \sigma_{Z}\right)$ in place of $(X, \rho)$ and with the property that for $f, f^{\prime} \in C_{\rho}, g \in C_{\rho \otimes \tau_{Z}}$, and $h \in C_{\rho \otimes \sigma_{Z}}$ we have

$$
\begin{aligned}
\left\langle D_{\rho}^{Z} f, g\right\rangle=(-1)^{p}\left\langle f, \theta E^{Z} g\right\rangle, & \left\langle E^{Z} f, h\right\rangle=(-1)^{p}\left\langle f, \theta D_{\rho \otimes \sigma_{Z}}^{Z} h\right\rangle, \\
\left\langle L_{\rho}^{Z} f, f^{\prime}\right\rangle=\left\langle f, L_{\rho}^{Z} f^{\prime}\right\rangle, & \left\langle M_{\rho}^{Z} f, f^{\prime}\right\rangle=\left\langle f, M_{\rho}^{Z} f^{\prime}\right\rangle, \\
\left\langle L_{\rho}^{Z} f, f\right\rangle \geq 0, & \left\langle M_{\rho}^{Z} f, f\right\rangle \geq 0
\end{aligned}
$$

under suitable convergence conditions (see below).
The theory of differential operators in this section was developed in [S80], [S81b], [S84a], [S84b], and [S86]. In particular, the last theorem was obtained in [S84a, Theorem 11.5, Corollary 11.8] and [S84b, p.486]. An equivalent result formulated on $G_{v}$ was given in [S90, Theorem 4.3, Corollary 4.4]. See also [S90, Lemma 4.2] for the definition of the inner product on $Z \otimes X$. The operators on $G_{v}$ correspond to $D_{\rho}^{Z}$ and $E^{Z}$, but the choice of $\Xi$ in $[\mathrm{S} 90,(7.7)]$ is the same as that of (12.4b) only for Types AB and CB; for Types AT and CT the latter is twice the former, and so the present $E^{Z}$ is $E$ of $[S 90,(7.9)]$ times a positive rational number, which may or may not be 1 . Thus, in order to have the first two equalities of the above theorem, we have to take this constant into account, though it does not play any essential role in practically all applications.

All these were formulated for functions on $\Gamma \backslash \mathcal{H}$ or on $\Gamma \backslash G_{\mathrm{a}}$ of compact support. To state a sufficient convergence condition in a more general case, we first note that given $f \in C_{\rho}(\Gamma)$, the function $\tilde{f}$ on $G_{\mathbf{a}}$ defined by $\tilde{f}(g)=\rho\left(M(g, \mathbf{o})^{-1}\right) f(g \mathbf{o})$ for $g \in G_{\mathbf{a}}$ is left $\Gamma$-invariant, where $\mathbf{o}=\left(\mathbf{o}_{v}\right)$ is a point of $\mathcal{H}$ such that $\left\{\alpha \in G_{v} \mid \alpha \mathbf{o}_{v}\right.$ $\left.=\mathbf{o}_{v}\right\}$ is the standard maximal compact subgroup of $G_{v}$ with which the results of [S90] are formulated. To be explicit, we take $\mathbf{o}_{v}=0$ for Types AB and CB, and $\mathbf{o}_{v}=i 1_{n}$ for Types AT and CT. Let $\mathfrak{g}_{v}$ be the Lie algebra of $G_{v}$ and $\mathfrak{g}_{v}=$ $\mathfrak{k}_{v}+\mathfrak{p}_{v}$ its Cartan decomposition. Given $\{\rho, X\},\left\{\rho^{\prime}, X^{\prime}\right\}, f \in C_{\rho}(\Gamma), h \in C_{\rho^{\prime}}(\Gamma)$, and a positive integer $p$, we say that ( $f, h$ ) is an integrable pair of type $(p, v)$, if $\psi\left(Y_{1} \cdots Y_{\mu} \tilde{f}\right) \psi^{\prime}\left(Y_{1}^{\prime} \cdots Y_{\nu}^{\prime} \tilde{h}\right)$ belongs to $L^{1}\left(\Gamma \backslash G_{\mathbf{a}}\right)$ for every $\psi \in S_{1}(X), \psi^{\prime} \in$ $S_{1}\left(X^{\prime}\right)$, and every $Y_{i}, Y_{j}^{\prime} \in \mathfrak{p}_{v}$ with $\mu \geq 0$ and $\nu \geq 0$ such that $\mu+\nu=p$ or $\mu+\nu=p-1$. Now the first resp. second formula of Theorem 12.15 is valid if
(12.36) $(f, g)$ resp. $(f, h)$ is an integrable pair of type $(p, v)$.

The reason for this is explained in [S94b, p.173]. As noted in [S90, p.257], the formulas are valid if $\tilde{f}, \tilde{g}, \tilde{h}$ are $C^{\infty}$ vectors in $L^{2}\left(\Gamma \backslash G_{\mathbf{a}}\right)$. Another sufficient condition is that all the holomorphic and anti-holomorphic derivatives of $f$ are rapidly decreasing and $g, h$ are slowly increasing at the cusps of $G$. Sufficient conditions for the last four formulas of Theorem 12.15 can be given in a similar manner or in the style of (12.36), since they are straightforward consequences of the first two formulas.

The relationship between the formulation on $H_{v}$ and that on $G_{v}$ is explained in [S90, §7] and [S94b, p.150]. See also Section A8. of the present book.
12.16. Corollary. Put $L_{\rho . v}=-\theta D_{\rho \otimes \sigma_{c}^{1} \cdot v} E_{v}$ for each $v \in \mathbf{a}^{\prime}$ ard put also $\Lambda=\sum_{v \in \mathbf{a}^{\prime}} c_{v} L_{\rho, v}$ with some fixed positive real numbers $c_{v}$, where $\mathbf{a}^{\prime}$ is defined by (12.5). Suppose $f \in C_{\rho}$ and $\left(f, E_{v} f\right)$ is an integrable pair of type $(1, v)$ for every $v \in \mathbf{a}^{\prime}$. Then $f$ is holomorphic if and only if $\Lambda f=0$.

Proof. If $f$ is holomorphic, then $E_{v} f=0$, and so $\Lambda f=0$. Conversely, if $\Lambda f=$ 0 , then by the second equality of Theorem 12.15 we have $\sum_{r \in \mathbf{a}^{\prime}} c_{r}\left\langle E_{r} \cdot f, E_{r} f\right\rangle=$ $\langle f, \Lambda f\rangle=0$ under the given integrability condition. Thus $E_{r} f=0$ for every $v \in \mathbf{a}^{\prime}$, and hence $f$ is holomorphic.
12.17. Let us now show that in the one-dimensional case the operators of this section have simple expressions. Thus take $G_{\mathbf{a}}=S L_{2}(\mathbf{R}), H=\{z \in \mathbf{C} \mid \operatorname{Im}(z)>$ $0\}, \mathfrak{K}_{0}=\mathbf{C}^{\times}$, and $T=\mathbf{C}$. Then for $z=x+i y \in T$ we have $\xi(z)=\eta(z)=2 y$. We first define operators $\varepsilon, \delta_{k}$, and $L_{k}$ on $H$ formally by

$$
\begin{align*}
& \varepsilon f=4 y^{2} \partial f / \partial \bar{z}, \quad \delta_{k} f=y^{-k}(\partial f / \partial z)\left(y^{k} f\right)  \tag{12.37}\\
& L_{k}=-\delta_{k-2} \varepsilon=-y^{2}\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)+2 i k y \frac{\partial}{\partial \bar{z}} \tag{12.38}
\end{align*}
$$

and define also $\delta_{k}^{p}$ for $0 \leq p \in \mathbf{Z}$ inductively by

$$
\begin{equation*}
\delta_{k}^{p+1}=\delta_{k+2 p} \delta_{k}^{p}, \quad \delta_{k}^{1}=\delta_{k}, \quad \delta_{k}^{0}=1 \tag{12.39}
\end{equation*}
$$

We now consider $\{\rho, \mathbf{C}\}$ with $\rho(a)=a^{k}$ with $k \in \mathbf{Z}$. We can put $S_{p}(T)=\mathbf{C}$ by identifying an element $h$ of $S_{p}(T)$ with $h(1)$. Then, dropping the subscript $v$, we have $D f=\partial f / \partial z, \bar{D} f=\partial f / \partial \bar{z}$, and (12.18) becomes $D_{\rho} f=y^{-k}(\partial f / \partial z)\left(y^{k} f\right)$, which is exactly $\delta_{k} f$ of (12.37). Similarly $E f=4 y^{2} \partial f / \partial \bar{z}=\varepsilon f$ with the above $\varepsilon$. Since $\left(\rho \otimes \tau^{p}\right)(a)=a^{k+2 p}$, from Proposition 12.10 (1) we see that $D_{\rho}^{p}$ coincides with the above $\delta_{k}^{p}$; thus Proposition 10.2 (2) specialized to the present case gives

$$
\begin{equation*}
\left(\delta_{k}^{p} f\right)\left\|_{k+2 p} \alpha=\delta_{k}^{p}\left(f \|_{k} \alpha\right), \quad(\varepsilon f)\right\|_{k-2} \alpha=\varepsilon\left(f \|_{k} \alpha\right) \quad\left(\alpha \in S L_{2}(\mathbf{R})\right) \tag{12.40}
\end{equation*}
$$

Also we have $Z=S_{p}(T)=\mathbf{C}$, and $L_{\rho}^{Z}$ of (12.31) for $Z=S_{1}(T)$ is exactly $L_{k}$ of (12.38), which has the property $\left(L_{k} f\right) \|_{k} \alpha=L_{k}\left(f \|_{k} \alpha\right)$ for every $\alpha \in S L_{2}(\mathbf{R})$.

## 13. Nearly holomorphic functions

We start with an easy lemma.
13.1. Lemma. Let $X_{1}, \ldots, X_{n}$ be mutually commutative $C^{\infty}$ vector fields on an $N$-dimensional $C^{\infty}$ manifold $U$ and $r_{1}, \ldots, r_{n}$ be elements of $C^{\infty}(U)$ such that the $n \times n$-matrix $\left(X_{j} r_{k}\right)$ is everywhere invertible; we assume that $n \leq N$. (Here the $X_{k}$ and the $r_{k}$ are complex-valued.) Define vector fields $Y_{1}, \ldots, Y_{n}$ by $Y_{i}=\sum_{j=1}^{n} b_{i j} X_{j}$ with the functions $b_{i j}$ determined by $\sum_{j=1}^{n} b_{i j} X_{j} r_{k}=\delta_{i k}$. Then the $Y_{i}$ are mutually commutative.

Proof. Since $X_{i} X_{j}=X_{j} X_{i}$, we can easily verify that $Y_{p} Y_{q}-Y_{q} Y_{p}=\sum_{j} c_{j}^{p q} X_{j}$ with $c_{j}^{p q} \in C^{\infty}(U)$. Now $Y_{p} Y_{q} r_{k}=Y_{p} \sum_{j=1}^{n} b_{q j} X_{j} r_{k}=Y_{p} \delta_{q k}=0$ for every $p$ and $q$, and hence $\sum_{j} c_{j}^{p q} X_{j} r_{k}=0$ for every $k$. Since $\operatorname{det}\left(X_{j} r_{k}\right) \neq 0$, we have $c_{j}^{p q}=0$, so that $Y_{p} Y_{q}-Y_{q} Y_{p}=0$ as expected.
13.2. We now consider an $n$-dimensional complex manifold $W$, and take $n$ elements $r_{1}, \ldots, r_{n}$ of $C^{\infty}(W)$ with the following property:
(13.1) Every point of $W$ has a small neighborhood $U$ on which there exist local complex coordinate functions $z_{1}, \ldots, z_{n}$ such that $\left(\partial r_{k} / \partial \bar{z}_{j}\right)_{j, k=1}^{n}$ is invertible everywhere on $U$.
Then we can define $2 n$ vector fields $\partial / \partial r_{k}$ and $\partial / \partial \bar{r}_{k}$ for $1 \leq k \leq n$ by the relations

$$
\begin{equation*}
\partial / \partial \bar{z}_{j}=\sum_{k=1}^{n}\left(\partial r_{k} / \partial \bar{z}_{j}\right) \partial / \partial r_{k}, \quad \partial / \partial z_{j}=\sum_{k=1}^{n}\left(\partial \bar{r}_{k} / \partial z_{j}\right) \partial / \partial \bar{r}_{k} . \tag{13.2}
\end{equation*}
$$

It can easily be seen that these vector fields do not depend on the choice of local coordinates, and so they are meaningful on the whole $W$. By the above lemma each set of $n$ vector fields $\partial / \partial r_{1}, \ldots, \partial / \partial r_{n}$ or $\partial / \partial \bar{r}_{1}, \ldots, \partial / \partial \bar{r}_{n}$ are mutually commutative; however, $\partial / \partial \bar{r}_{i}$ and $\partial / \partial \bar{r}_{j}$ do not necessarily commute, as can easily be seen by a counterexample, even when $n=1$. In view of the commutativity, we can naturally speak of $\partial^{m} / \partial r_{i_{1}} \cdots \partial r_{i_{m}}$ and $\partial^{m} / \partial \bar{r}_{i_{1}} \cdots \partial \bar{r}_{i_{m}}$ with no ambiguity.

The reason for choosing the symbols $\partial / \partial r_{i}$ and $\partial / \partial \bar{r}_{i}$ for these vector fields will be explained by the first assertion of the following lemma.
13.3. Lemma. (1) Let $g(z, w)$ be a $C^{\infty}$ function of $(z, w) \in W \times V$ with an open subset $V$ of $\mathbf{C}^{n}$, holomorphic in $w$, and let $h(z, w)=\partial^{e} g / \partial w_{i_{1}} \cdots \partial w_{i_{e}}$. If $g$ is holomorphic in $z$, then $h(z, r(z))=\left(\partial^{e} / \partial r_{i_{1}} \cdots \partial r_{i_{e}}\right) g(z, r(z))$ whenever $g(z, r(z))$ is meaningful. Similarly, if $g(z, w)$ is antiholomorphic in $z \in W$, then $h(z, \bar{r}(z))=\left(\partial^{e} / \partial \bar{r}_{i_{1}} \cdots \partial \bar{r}_{i_{e}}\right) g(z, \bar{r}(z))$.
(2) The $r_{i}$ are algebraically independent over the field of all meromorphic functions on $W$.
(3) $A C^{\infty}$ function $f$ on $W$ is a polynomial in $r_{1}, \ldots, r_{n}$ of degree $<e$ with holomorphic functions as coefficients if and only if all the derivatives of the form $\partial^{e} f / \partial r_{i_{1}} \cdots \partial r_{i_{e}}$ are 0 .

Proof. Since $\partial / \partial r_{i}$ (resp. $\partial / \partial \bar{r}_{i}$ ) annihilates holomorphic (resp. antiholomorphic) functions, assertion (1) can be verified easily. To prove (2), let $X=$ $\left(X_{1}, \ldots, X_{n}\right)$ be a set of indeterminates, and $Q(X)$ a polynomial of degree $\leq m$ whose coefficients are holomorphic functions on $W$. We prove by induction on $m$ that $Q=0$ if $Q(r)=0$. This is trivial if $m=0$. Suppose $Q(r)=0$ with $m>0$. Then by (1), $\left(\partial Q / \partial X_{i}\right)(r)=\left(\partial / \partial r_{i}\right) Q(r)=0$, and hence $\partial Q / \partial X_{i}=0$ by the induction assumption. Thus $Q$ is a constant, and must be 0 . This proves (2). We prove the "if"-part of (3) by induction on $e$. The case $e=1$ is obvious. Suppose all the derivatives of the form $\partial^{e} f / \partial r_{i_{1}} \cdots \partial r_{i_{e}}$ are 0 with some $e>1$. Then by the induction assumption for $e-1$, we can find, for each $i$, a polynomial $Q_{i}(X)$ of degree $<e-1$ whose coefficients are holomorphic functions on $W$ such that $\partial f / \partial r_{i}=Q_{i}(r)$. Put $Q_{i j}=\partial Q_{i} / \partial X_{j}$. By (1), we have $Q_{i j}(r)=\partial^{2} f / \partial r_{i} \partial r_{j}=Q_{j i}(r)$, and hence $Q_{i j}=Q_{j i}$ by (2). Therefore we can find a polynomial $P(X)$ of degree $<e$ whose coefficients are holomorphic functions on $W$ such that $\partial P / \partial X_{i}=Q_{i}$ for every $i$. (Indeed, for $x=\left(x_{i}\right) \in \mathbf{R}^{n}$ and $a \in \mathbf{R}^{n}$ define $P(a)=\int_{c} \sum_{i=1}^{n} Q_{i}(x) d x_{i}$ with any path $c$ connecting 0 to $a$. Then $\partial P / \partial x_{i}=Q_{i}(x)$. Since $P(x)=\int_{0}^{1} \sum_{i=1}^{n} Q_{i}(t x) x_{i} d t$, we easily see that $P$ is a polynomial of degree $<e$ with holomorphic functions on $W$ as coefficients.) Then $\left(\partial / \partial r_{i}\right)(f-P(r))=0$, and hence $f-P(r)$ is holomorphic. This completes our induction. The "only if"-part of (3) follows immediately from (1).
13.4. Let $W$ be a complex Kähler manifold of complex dimension $N$ with a fundamental 2-form $\Omega$; let $\Omega=i \sum_{p, q=1}^{N} h_{p q} d z_{q} \wedge d \bar{z}_{p}$ with local complex coordinate functions $z_{1}, \ldots, z_{N}$ in a coordinate neighborhood $U$. We can then define $N$ vector fields $R_{1}, \ldots, R_{N}$ on $U$ by the relations

$$
\begin{equation*}
\partial / \partial \bar{z}_{p}=\sum_{q=1}^{N} h_{p q} R_{q} \quad(1 \leq p \leq N) \tag{13.3}
\end{equation*}
$$

It is well-known that $\Omega$, for a sufficiently small $U$, can be given in the form $\Omega=$ $i \sum_{p, q=1}^{N} \partial^{2} \varphi / \partial z_{q} \partial \bar{z}_{p} \cdot d z_{q} \wedge d \bar{z}_{p}$ with a real-valued function $\varphi$ on $U$. Define $N$ functions $r_{q}$ on $U$ by $r_{q}=\partial \varphi / \partial z_{q}$. Then $h_{p q}=\partial r_{q} / \partial \bar{z}_{p}$, and therefore we can apply the principle of $\S 13.2$ to the present case to find that $R_{q}=\partial / \partial r_{q}$.

Now the $R_{p}$ may depend on the choice of the $z_{p}$. However, if $w_{1}, \ldots, w_{N}$ are coordinate functions in a coordinate neighborhood $U^{\prime}$ and vector fields $S_{1}, \ldots, S_{N}$ are defined relative to the $w_{p}$ in the same manner, then $R_{p}=\sum_{q} \partial z_{p} / \partial w_{q} \cdot S_{q}$ on $U \cap U^{\prime}$. Therefore if we denote by $\mathcal{N}^{e-1}(U)$ the set of all elements of $C^{\infty}(U)$
annihilated by $X_{\nu_{1}} \cdots X_{\nu_{e}}$ for all $\left(\nu_{1}, \ldots, \nu_{e}\right) \in\{1, \ldots, N\}^{e}$, this is well-defined independently of the choice of the $z_{p}$. Then, for every open subset $V$ of $W$ let $\mathcal{N}^{e}(V)$ denote the set of all $f \in C^{\infty}(V)$ such that the restriction of $f$ to any coordinate neighborhood $U$ belongs to $\mathcal{N}^{e}(U)$, and let $\mathcal{N}(V)=\bigcup_{e=0}^{\infty} \mathcal{N}^{e}(V)$. We call an element of $\mathcal{N}(V)$ a nearly holomorphic function on $V$ relative to $\Omega$. If the $r_{p}$ are defined on $U$ as above, then Lemma 13.3 (3) shows that $\mathcal{N}^{e}(U)$ consists of all the polynomials in the $r_{p}$ whose coefficients are holomorphic functions on $U$. Clearly $\mathcal{N}^{0}(U)$ is the set of all holomorhic functions on $U$. In our applications $W$ is a hermitian symmetric space, and we can define the functions $r_{p}$ on the whole $W$, so that the last statement is applicable to $\mathcal{N}^{e}(W)$.
13.5. We now consider the space $\mathcal{H}$ of (12.4a). We first define a matrix-valued function $r_{v}$ and a scalar-valued function $\delta_{v}$ on $H_{v}$ as follows:

$$
\begin{align*}
& r_{v}(z)=-\xi(z)^{-1} \bar{z}=-\bar{z} \cdot{ }^{t} \eta(z)^{-1} \quad \text { (Types AB, CB), }  \tag{13.4a}\\
& r_{v}(z)=\left({ }^{t} z-\bar{z}\right)^{-1} \quad \text { (Types AT, CT), }  \tag{13.4b}\\
& \delta_{v}(z)= \begin{cases}\operatorname{det}\left[2^{-1} \eta(z)\right] & (\text { Types AT, CT) }, \\
\operatorname{det}[\eta(z)] & (\text { Types AB, CB). }\end{cases} \tag{13.5}
\end{align*}
$$

The function $\delta_{v}$ is the same as what was defined in (3.21). We then put $r(z)=$ $\left(r_{v}\left(z_{v}\right)\right)_{v \in \mathbf{a}}$ and $\delta(z)=\left(\delta_{v}\left(z_{v}\right)\right)_{v \in \mathbf{a}}$ for $z=\left(z_{v}\right)_{v \in \mathbf{a}} \in \mathcal{H}$; we do not define $r_{v}$ if $G_{v}$ is compact. We fix one $v \in \mathbf{a}$ and consider the behavior of a function on $\mathcal{H}$ only with respect to $z_{v}$. For simplicity let us drop temporarily the subscript $v$ from the objects $T_{v}, D_{v}, \bar{D}_{v}, \xi_{v}$, etc. Then for $u \in T$ we have

$$
\begin{align*}
& \eta^{-1}(D \eta)(u)={ }^{t} r u  \tag{13.6a}\\
& \xi^{-1}(D \xi)(u)=r \cdot{ }^{t} u  \tag{13.6b}\\
& (D r)(u)=-r \cdot{ }^{t} u r  \tag{13.6c}\\
& (\bar{D} r)(u)=-\xi^{-1} u \cdot{ }^{t} \eta^{-1} .
\end{align*}
$$

The first two formulas can be verified in a straightforward way. As for (13.6d) for Type AB, for example, we have $\xi r=-\bar{z}$, and so $(\bar{D} \xi)(u) r+\xi(\bar{D} r)(u)=-u$. Thus $(\bar{D} r)(u)=-\xi^{-1}(u+(\bar{D} \xi)(u) r)=-\xi^{-1}\left(u-u \cdot{ }^{t} z r\right)=-\xi^{-1} u\left({ }^{t} \eta+{ }^{t} z \bar{z}\right) \cdot{ }^{t} \eta^{-1}=$ $-\xi^{-1} u \cdot{ }^{t} \eta^{-1}$. All the remaining cases can be verified in the same fashion.
13.6. Lemma. Let $\zeta$ and $\psi_{Z}$ be as in Theorem 12.13; let $s \in \mathbf{C}$. Then

$$
\begin{equation*}
\zeta(\mathcal{D}) \delta(z)^{s}=\psi_{Z}(s) \delta(z)^{s} \zeta(r(z)) \tag{13.7}
\end{equation*}
$$

Proof. For Types AB and CB we have $\delta(z)=\operatorname{det}\left(1_{n}-z^{*} z\right)$. Taking $\left(-z^{*}, 1\right)$ to be ( $c, d$ ) in Theorem 12.13, we obtain the desired formula. The other cases can be handled in the same manner.
13.7. We now note three basic formulas:

$$
\begin{align*}
& D\left(\delta^{s}\right)(u)=s \cdot \delta^{s} \operatorname{tr}\left({ }^{t} r u\right),  \tag{13.8a}\\
& (D \log \delta)(u)=\operatorname{tr}\left({ }^{t} r u\right),  \tag{13.8b}\\
& (\bar{D} D \log \delta)(u, v)=-\operatorname{tr}\left({ }^{t} u \xi^{-1} v \cdot{ }^{t} \eta^{-1}\right) \quad(u, v \in T) \tag{13.8c}
\end{align*}
$$

Indeed, take $Z=S_{1}(T)$ and $\zeta(x)=\operatorname{tr}\left({ }^{t} u x\right)$ with $u \in T$ in Lemma 13.6. Then $\zeta\left(\mathcal{D}_{v}\right) f=(D f)(u)$, and so we obtain (13.8a). The second equality follows immediately from the first one. Combining it with (13.6d), we obtain (13.8c).

Since $\xi$ and $\eta$ are hermitian and positive definite, we see that $H$ is a Kähler manifold with $i \bar{\partial} \partial \log \delta$ as its fundamental 2 -form. For the product space $\mathcal{H}=$ $\prod_{v \in \mathbf{a}} H_{v}$, we have to take $i \sum_{v \in \mathbf{a}} \bar{\partial} \partial \log \delta_{v}\left(z_{v}\right)$. Then (13.8b) shows that the entries of the functions $\left(r_{v}\right)_{v \in \mathbf{a}}$ are exactly the $r_{p}$ discussed in §13.4.
13.8. Take an $\mathbf{R}$-rational basis $\left\{\varepsilon_{\nu}\right\}_{\nu \in N}$ of $T_{v}$ over $\mathbf{C}$ and put $z_{v}=\sum_{\nu \in N} z_{v \nu} \varepsilon_{\nu}$ and $r_{v}=\sum_{\nu \in N} r_{v \nu} \varepsilon_{\nu}$. Following the general principle of $\S 13.2$, define vector fields $\partial / \partial r_{v \nu}$ and $\partial / \partial \bar{r}_{v \nu}$ by

$$
\begin{equation*}
\partial / \partial z_{v \mu}=\sum_{\nu \in N}\left(\partial \bar{r}_{v \nu} / \partial z_{v \mu}\right) \partial / \partial \bar{r}_{v \nu}, \quad \partial / \partial \bar{z}_{v \mu}=\sum_{\nu \in N}\left(\partial r_{v \nu} / \partial \bar{z}_{v \mu}\right) \partial / \partial r_{v \nu} \tag{13.9}
\end{equation*}
$$

These are well-defined in view of (13.6d). Now, for $u=\sum_{\nu \in N} u_{v \nu} \varepsilon_{\nu} \in T_{v}$, we have

$$
\begin{equation*}
\left(C_{v} f\right)(u)=-\sum_{\nu \in N} u_{v \nu} \partial f / \partial \bar{r}_{v \nu}, \quad\left(E_{v} f\right)(u)=-\sum_{\nu \in N} u_{v \nu} \partial f / \partial r_{v \nu} \tag{13.10}
\end{equation*}
$$

Indeed, $\left(D_{v} f\right)(u)=\sum_{\nu \in N} u_{v \nu} \partial f / \partial z_{v \nu}=\sum_{\nu \in N}\left(D_{v} \bar{r}_{v \nu}\right)(u) \partial f / \partial \bar{r}_{v \nu}$, and therefore $\left(D_{v} f\right)\left({ }^{t} \xi_{v} u \eta_{v}\right)=\sum_{\nu \in N}\left(D_{v} \bar{r}_{v \nu}\right)\left({ }^{t} \xi_{v} u \eta_{v}\right) \partial f / \partial \bar{r}_{v \nu}$, which together with (13.6d) and (12.12b) proves the first equality of (13.10); the second one can be proved similarly.

Since the vector fields $\partial / \partial \bar{r}_{v \nu}$ mutually commute, from (13.10) we see that the values of $C_{v}^{e} f$ are symmetric elements of $M l_{e}\left(T_{v}, X\right)$; the same is true for $E_{v}^{e} f$ because of the commutativity of the $\partial / \partial r_{v \nu}$.

Taking the basis $\left\{\varepsilon_{\nu}^{\prime}\right\}_{\nu \in N}$ of $T_{v}$ dual to $\left\{\varepsilon_{\nu}\right\}_{\nu \in N}$, we define symbols $\partial / \partial r_{v}$ and $\partial / \partial \bar{r}_{v}$ by

$$
\begin{equation*}
\partial / \partial r_{v}=\sum_{\nu \in N} \varepsilon_{\nu}^{\prime} \partial / \partial r_{v \nu}, \quad \partial / \partial \bar{r}_{v}=\sum_{\nu \in N} \varepsilon_{\nu}^{\prime} \partial / \partial \bar{r}_{v \nu} . \tag{13.11}
\end{equation*}
$$

These are independent of the choice of bases of $T_{v}$. Given $g \in S_{p}\left(T_{v}\right)$, we define $g\left(\partial / \partial r_{v}\right)$ by
(13.12) $g\left(\partial / \partial r_{v}\right)=g_{*}\left(\partial / \partial r_{v}, \ldots, \partial / \partial r_{v}\right)=\sum g_{*}\left(\varepsilon_{\nu_{1}}^{\prime}, \ldots, \varepsilon_{\nu_{\nu}}^{\prime}\right) \partial^{p} / \partial r_{v \nu_{1}} \cdots \partial r_{v \nu_{p}}$, where $\left(\nu_{1}, \ldots, \nu_{p}\right)$ runs over $N^{p}$, and define $g\left(\partial / \partial \bar{r}_{v}\right)$ similarly. Then we have

$$
\begin{equation*}
g\left(\partial / \partial \bar{r}_{v}\right) f=(-1)^{p}\left[g, C_{v}^{p} f\right], \quad g\left(\partial / \partial r_{v}\right) f=(-1)^{p}\left[g, E_{v}^{p} f\right] . \tag{13.13}
\end{equation*}
$$

These follow from (12.8) and (13.10) immediately.
The notation being the same as in (12.20) and (12.23), we have

$$
\begin{gather*}
\left(D_{\rho}^{Z} f\right)(\zeta)=(-1)^{p} \rho(\Xi)^{-1} \zeta^{\prime}\left(\partial / \partial \bar{r}_{v}\right)(\rho(\Xi) f)  \tag{13.14a}\\
\left(E^{Z} f\right)(\zeta)=(-1)^{p} \zeta\left(\partial / \partial r_{v}\right) f
\end{gather*}
$$

for every $\zeta \in Z$, where $\zeta^{\prime}=\sigma_{v}^{p}(\Xi) \zeta$. Indeed, combining (12.23) with (12.17) and (12.9), we obtain

$$
\left(D_{\rho}^{Z} f\right)(\zeta)=\left[\zeta,\left(\rho \otimes \tau_{v}^{p}\right)(\Xi)^{-1} C_{v}^{p}(\rho(\Xi) f)\right]=\rho(\Xi)^{-1}\left[\sigma_{v}^{p}(\Xi) \zeta, C_{v}^{p}(\rho(\Xi) f)\right]
$$

Applying (13.13) to the last quantity, we obtain (13.14a). Formula (13.14b) follows directly from (12.23) and (13.13).
13.9. Lemma. Let $\rho(a, b)=\operatorname{det}(b)^{k}$ and $f(z)=\left(\prod_{v \in \mathbf{a}} \delta_{v}\left(z_{v}\right)^{s}\right) \|_{k} \alpha$ with $k \in$ $\mathbf{Z}^{\mathbf{a}}, \alpha \in G_{\mathbf{a}}$, and $s \in \mathbf{C}$; let $Z$ be an irreducible subspace of $S_{p}\left(T_{v}\right)$ and $\psi_{Z}$ be as in Theorem 12.13. Then for every $\zeta \in Z$ we have

$$
\begin{aligned}
& \left(D_{\rho}^{Z} f\right)(\zeta)=\left\{\begin{array}{lr}
i^{p} \psi_{Z}\left(-k_{v}-s\right) \zeta\left(\left(\xi^{-1} \lambda_{\alpha}^{*} \cdot{ }^{t} \mu_{\alpha}^{-1}\right)_{v}\right) f & \text { (Types AT, CT) } \\
\psi_{Z}\left(-k_{v}-s\right) \zeta\left(\left(\xi^{-1} \lambda_{\alpha}^{*} \cdot \overline{\alpha z} \cdot{ }^{t} \mu_{\alpha}^{-1}\right)_{v}\right) f & \text { (Types AB, CB) }
\end{array}\right. \\
& \left(E^{Z} f\right)(\zeta)= \begin{cases}(-i)^{p} \psi_{Z}(-s) \zeta\left(\left({ }^{( } \lambda_{\alpha} \widehat{\mu}_{\alpha} \eta\right)_{v}\right) f & \text { (Types AT, CT) } \\
\psi_{Z}(-s) \zeta\left(\left({ }^{t}(\alpha z) \lambda_{\alpha} \widehat{\mu}_{\alpha} \eta\right)_{v}\right) f & \text { (Types AB, CB) }\end{cases}
\end{aligned}
$$

Proof. Clearly it is sufficient to consider only $\delta_{v}\left(z_{v}\right)^{s}$; so we take $\alpha \in G_{v}$ and drop the subscript $v$. Put $\alpha=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right]$ with $d$ of size $n$. For Types AT and CT we have $\bar{r}=\left(z^{*}-z\right)^{-1}$, and therefore

$$
\rho(\Xi) f=2^{n k} \delta^{k+s} j_{\alpha}^{-k-s}\left(\bar{j}_{\alpha}\right)^{-s}=2^{-n s} i^{n(k+s)}\left(\bar{j}_{\alpha}\right)^{-s} \operatorname{det}((c z+d) \bar{r})^{-k-s}
$$

Since $z=z^{*}-\bar{r}^{-1}$, we have $(c z+d) \bar{r}=\left(c z^{*}+d\right) \bar{r}-c$. Now $\zeta^{\prime}(\partial / \partial \bar{r}) g=0$ for every anti-holomorphic $g$, so that for every $\zeta^{\prime} \in Z$,

$$
\begin{equation*}
\zeta^{\prime}(\partial / \partial \bar{r})(\rho(\Xi) f)=2^{-n s} i^{n(k+s)}\left(\bar{j}_{\alpha}\right)^{-s} \zeta^{\prime}(\partial / \partial \bar{r}) \operatorname{det}(A \bar{r}+B)^{-k-s} \tag{*}
\end{equation*}
$$

where $(A, B)=\left(c z^{*}+d,-c\right)$. In view of Lemma 13.3 (1) we can apply Theorem 12.13 with $\partial / \partial \bar{r}$ as $\mathcal{D}$ there. (Notice that $A \cdot{ }^{t} B=B \cdot{ }^{t} A$ for Type CT.) Thus the quantity of $\left({ }^{*}\right)$ equals

$$
\begin{equation*}
\psi_{Z}(-k-s) \rho(\Xi) f \cdot \zeta^{\prime}\left({ }^{t} A \cdot{ }^{t}(A \bar{r}+B)^{-1}\right) \tag{**}
\end{equation*}
$$

Now we employ (13.14a) in the present setting. Then our task is to calculate $\zeta^{\prime}\left({ }^{t} A \cdot{ }^{t}(A \bar{r}+B)^{-1}\right)$ with $\zeta^{\prime}(u)=\zeta\left(\xi^{-1} u \cdot{ }^{t} \eta^{-1}\right)$. Since ${ }^{t} A \cdot{ }^{t}(A \bar{r}+B)^{-1}={ }^{t}\left(c z^{*}+\right.$ $d) \cdot{ }^{t}(c z+d)^{-1} \widehat{r}$, we obtain the formula for $\left(D_{\rho}^{Z} f\right)(\zeta)$ for Types AT and CT.

The formula for $\left(E^{Z} f\right)(\zeta)$ is simpler. Since $\bar{z}={ }^{t} z-r^{-1}$, we see that $\delta^{s} \|_{k} \alpha$ is a holomorphic factor times $\operatorname{det}\left[\left(\bar{c} \cdot{ }^{t} z+\bar{d}\right) r-\bar{c}\right]^{-s}$. Applying $\zeta(\partial / \partial r)$ to this, from (13.14b) and Theorem 12.13 we obtain the desired formula.

The argument for Types AB and CB is similar but requires modifications. We have $1-z^{*} \bar{r}=\left(\eta+z^{*} z\right) \eta^{-1}=\eta^{-1}$, and so $(c z+d) \eta^{-1}=c z \eta^{-1}+d \eta^{-1}=$ $-c \bar{r}+d\left(1-z^{*} \bar{r}\right)=\left(-c-d z^{*}\right) \bar{r}+d$, and so

$$
\rho(\Xi) f=\delta^{k+s} j_{\alpha}^{-k-s}\left(\bar{j}_{\alpha}\right)^{-s}=\left(\bar{j}_{\alpha}\right)^{-s} \operatorname{det}(A \bar{r}+B)^{-k-s}
$$

with $(A, B)=\left(-c-d z^{*}, d\right)$. By Theorem 12.13, $\zeta^{\prime}(\partial / \partial \bar{r})(\rho(\Xi) f)$ equals $\left({ }^{* *}\right)$. From (3.14) we see that $A=-(\alpha z)^{*} \overline{\lambda_{\alpha}(z)}$, and hence $-\xi^{-1} \cdot{ }^{t} A \cdot{ }^{t}(A \bar{r}+B)^{-1} \cdot{ }^{t} \eta^{-1}=$ $\xi^{-1} \lambda_{\alpha}^{*} \cdot(\overline{\alpha z}) \cdot{ }^{t} \mu_{\alpha}^{-1}$, which gives the desired formula for $\left(D_{\rho}^{Z} f\right)(\zeta)$. As for $\left(E^{Z} f\right)(\zeta)$, we observe that $1-{ }^{t} z r={ }^{t} \eta^{-1}$, and reduce the problem to $\zeta(\partial / \partial r) \operatorname{det}[(-\bar{c}-\bar{d}$. $\left.\left.{ }^{t} z\right) r+d\right]^{-s}$. Then the formula can be obtained in a similar fashion.

In the above proof for Types AB and CB we obtained $c+d z^{*}$. This is not a factor of automorphy, but may be called a quasi-factor of automorphy for the following reason. Put $\kappa_{\alpha}(z)=\bar{c}+\bar{d} \cdot{ }^{t} z$. Since $\overline{\kappa_{\alpha}(z)}$ is the lower left $n \times m$-block of $\alpha B(z)$ with $B(z)$ of (3.10), we can easily verify that

$$
\begin{equation*}
\kappa_{\alpha}(z)={ }^{t}(\alpha z) \lambda_{\alpha}(z), \quad \kappa_{\alpha \beta}(z)=\kappa_{\alpha}(\beta z) \lambda_{\beta}(z) \tag{13.15}
\end{equation*}
$$

Also we can put $\zeta\left(\left(\xi^{-1} \lambda_{\alpha}^{*} \cdot \overline{\alpha z} \cdot{ }^{t} \mu_{\alpha}^{-1}\right)_{v}\right)=\zeta\left(\left(\xi^{-1} \kappa_{\alpha}^{*} \cdot{ }^{t} \mu_{\alpha}^{-1}\right)_{v}\right)$. See $[\mathrm{S} 86, \S \S 5,6,(5.5)$ and (5.6) in particular] for the formulas essentially of the same nature given in a somewhat different manner.
13.10. Lemma. The difference $r_{v}\left(\alpha_{v} z_{v}\right)-\lambda\left(\alpha_{v}, z_{v}\right) r_{v}\left(z_{v}\right) \cdot{ }^{t} \mu\left(\alpha_{v}, z_{v}\right)$ is holomorphic in $z$ for every $\alpha \in G_{\mathbf{a}}$.

Proof. Dropping the subscript $v$ for simplicity and denoting by $1_{T}$ the identity map $T \rightarrow T$, we see from (13.10) that $E r=-1_{T}$. Define $\rho: K^{c} \rightarrow G L(T)$ by $\rho(a, b) u=a u \cdot{ }^{t} b$. Then $r \|_{\rho} \alpha=\lambda(\alpha, z)^{-1} r(\alpha z) \cdot{ }^{t} \mu(\alpha, z)^{-1}$. Clearly $[(\rho \otimes$ $\left.\left.\sigma^{1}\right)(a, b)\left(-1_{T}\right)\right]=-1_{T}$, and hence by Proposition 12.10 (2) we have $E\left(r \|_{\rho} \alpha\right)=$ $(E r) \|_{\rho \otimes \sigma^{1}} \alpha=E r$. Thus $E\left(r \|_{\rho} \alpha-r\right)=0$, so that $r \|_{\rho} \alpha-r$ is holomorphic as expected.
13.11. Given $p=\left(p_{v}\right)_{v \in \mathbf{a}} \in \mathbf{Z}^{\mathbf{a}}$ with $p_{v} \geq 0$ for every $v$ and a representation $\{\rho, X\}$ of $\mathfrak{K}_{0}$, we denote by $\mathcal{N}^{p}(\mathcal{H}, X)$ the set of all $f \in C^{\infty}(\mathcal{H}, X)$ which are polynomials in the components of $r=\left(r_{v}\right)_{v \in \mathbf{a}}$, of degree $\leq p_{v}$ in $r_{v}$, with holomorphic maps of $\mathcal{H}$ into $X$ as coefficients. We naturally take $p_{v}=0$ if $v \notin \mathbf{a}^{\prime}$. Then, (13.10) together with Lemma 13.3 (3) shows that $\mathcal{N}^{p}(\mathcal{H}, X)$ consists of all $f \in C^{\infty}(\mathcal{H}, X)$ such that $E_{v}^{p_{v}+1} f=0$ for every $v \in \mathbf{a}^{\prime}$. Moreover, by Lemma 13.3 (2) the components of $r_{v}$ for all $v$ are algebraically independent over the field of all meromorphic functions on $\mathcal{H}$. For example, if we view an element $f$ of $\mathcal{N}^{p}(\mathcal{H}, X)$ as a function of $z_{v}$ and suppress other variables $z_{v^{\prime}}$ for $v^{\prime} \in \mathbf{a}, \neq v$, then

$$
\begin{equation*}
f(z)=\sum_{i=0}^{p_{v}} g_{i}\left(z_{v}, r_{v}\left(z_{v}\right)\right) \tag{13.16}
\end{equation*}
$$

with a holomorphic map $g_{i}: H_{v} \rightarrow S_{i}\left(T_{v}, X\right)$ for each $i$, where $g_{i}\left(z_{v}, u\right)$ means the element $g_{i}\left(z_{v}\right)$ of $S_{i}\left(T_{v}, X\right)$ evaluated at $u \in T_{v}$. From Proposition 12.10 (2) we see that $\mathcal{N}^{p}(\mathcal{H}, X)$ is stable under the maps $f \mapsto f \circ \alpha$ and $f \mapsto f \|_{\rho} \alpha$ for every $\alpha \in G_{\mathbf{a}}$. The elements of $\bigcup_{p} \mathcal{N}^{p}(\mathcal{H}, X)$ are called ( $X$-valued) nearly holomorphic functions on $\mathcal{H}$, as defined at the beginning of this section. We note here
(13.17) With $g_{i}$ as in (13.16) we have $E_{v}^{p_{v}} f=(-1)^{p_{v}} p_{v}!g_{p_{v}}$, that is, $\left(E_{v}^{p_{v}} f\right)(u)=$ $(-1)^{p_{v}} p_{v}!g_{p_{v}}\left(z_{v}, u\right)$ for $u \in T_{v}$.
Indeed, $E_{v}^{p_{v}}$ kills $g_{i}\left(z_{v}, r_{v}\left(z_{v}\right)\right)$ for $i<p_{v}$. Now, by (13.13), $\left[h, E_{v}^{p_{v}} g_{p_{v}}\left(z_{v}, r_{v}\left(z_{v}\right)\right)\right]$ $=(-1)^{p_{v}} h\left(\partial / \partial r_{v}\right) g_{p_{v}}\left(z_{v}, r_{v}\left(z_{v}\right)\right)$ for every $h \in S_{p_{v}}\left(T_{v}\right)$. By (12.28) and Lemma 13.3 the last quantity is $(-1)^{p_{v}} p_{v}!\left[h, g_{p_{v}}\left(z_{v}, *\right)\right]$, which proves (13.17).

For a congruence subgroup $\Gamma$ of $G$ we denote by $\mathcal{N}_{\rho}^{p}(\Gamma)$ the subset of $C_{\rho}(\Gamma) \cap$ $\mathcal{N}^{p}(\mathcal{H}, X)$ consisting of the functions satisfying the cusp condition, which is required only when $G$ is isogenous to $S L_{2}(\mathbf{Q})$. (For the precise statement of the cusp condition see $\S 13.12$ below.) We then denote by $\mathcal{N}_{\rho}^{p}$ the union of $\mathcal{N}_{\rho}^{p}(\Gamma)$ for all $\Gamma$. Clearly $\mathcal{N}_{\rho}^{p} \|_{\rho} \alpha=\mathcal{N}_{\rho}^{p}$ for every $\alpha \in G$. Since $\mathcal{N}^{0}(\mathcal{H}, X)$ consists of all holomorphic maps of $\mathcal{H}$ to $X$, we see that $\mathcal{N}_{\rho}^{0}(\Gamma)=\mathcal{M}_{\rho}(\Gamma)$ and $\mathcal{N}_{\rho}^{0}=\mathcal{M}_{\rho}$ with the symbol $\mathcal{M}_{\rho}$ of Section 5.

In this book we consider almost exclusively nearly holomorphic functions on $\mathcal{H}$ of the above type, which form the most important class from the number-theoretical viewpoint. However, we can also determine such functions on hermitian symmetric spaces of compact type. For details, the reader is referred to [S87a].
13.12. Continuing the discussion of $\S 12.17$, let us now consider nearly holomorphic functions on the upper half plane $H$. In this case $r=(2 i y)^{-1}$, and so $\mathcal{N}^{p}(H, \mathbf{C})$ consists of the functions of the form $\sum_{\nu=0}^{p} y^{-\nu} g_{\nu}(z)$ with holomorphic functions $g_{\nu}$ on $H$, as we already mentioned in the introduction. Since $\partial r / \partial \bar{z}=-\left(4 y^{2}\right)^{-1}$, we have

$$
\begin{equation*}
\partial / \partial r=-4 y^{2} \partial / \partial \bar{z}, \quad \partial / \partial \bar{r}=-4 y^{2} \partial / \partial z \tag{13.18}
\end{equation*}
$$

We say that an element $f$ of $\mathcal{N}^{p}(H, \mathbf{C}) \cap C_{\rho}(\Gamma)$, for $\Gamma \subset S L_{2}(\mathbf{Q})$ and $\rho(x)=x^{k}$ with $k \in \mathbf{Z}$, satisfies the cusp condition if it satisfies
(13.18a) For every $\alpha \in S L_{2}(\mathbf{Z})$ we have

$$
\left(f \|_{k} \alpha\right)(z)=\sum_{\nu=0}^{p}(\pi y)^{-\nu} \sum_{n=0}^{\infty} c_{\alpha \nu n} \exp \left(2 \pi i n z / N_{\alpha}\right)
$$

with $c_{\alpha \nu n} \in \mathbf{C}$ and $0<N_{\alpha} \in \mathbf{Z}$.
13.13. Let $\{\rho, X\}$ be a representation of $\mathfrak{K}_{0}$. So far we considered $D_{\rho, v}^{p}$ and $D_{\rho}^{Z}$ for a fixed $v \in \mathbf{a}, 0 \leq p \in \mathbf{Z}$, and $Z \subset S_{p}\left(T_{v}\right)$. We now generalize this by considering the derivatives with respect to the variables on the whole $\mathcal{H}$. We put $T=\prod_{v \in \mathbf{a}} T_{v}$ and $T^{e}=\prod_{v \in \mathbf{a}} T_{v}^{e_{v}}$ for $e=\left(e_{v}\right)_{v \in \mathbf{a}} \in \mathbf{Z}^{\mathbf{a}}$ with $e_{v} \geq 0$; we take $e_{v}=0$ if $v \notin \mathbf{a}^{\prime}$. We then denote by $M l_{e}(T, X)$ the vector space of all $\mathbf{C}$-multilinear maps of $T^{e}$ (that is, C-linear on each single factor $T_{v}$ of $T^{e}$ ) into $X$, and by $S_{e}(T, X)$ the vector space of all polynomial maps of $T$ into $X$ homogeneous of degree $e_{v}$ in the variable on $T_{v}$ for each $v$. In particular we put $S_{e}(T)=S_{e}(T, \mathbf{C})$. Given $h \in S_{e}(T, X)$, there exists a unique element $h_{*}$ of $M l_{e}(T, X)$ which is symmetric on $T_{v}^{e_{v}}$ for every $v \in \mathbf{a}$, and such that

$$
h(y)=h_{*}(\overbrace{y_{v}, \ldots, y_{v}}^{e_{v}}, \overbrace{y_{w}, \ldots, y_{w}}^{e_{w}}, \ldots \ldots)
$$

for $y=\left(y_{v}\right)_{v \in \mathbf{a}} \in T$. This is an easy generalization of Lemma 12.4 (2). For $g \in S_{e}(T)$ and $g \in S_{e}(T, X)$ we define $[g, h] \in X$ by an obvious generalization of (12.8); also we define representations $\left\{\rho \otimes \tau^{e}, S_{e}(T, X)\right\}$ and $\left\{\rho \otimes \sigma^{e}, S_{e}(T, X)\right\}$ of $\mathfrak{K}_{0}$ by

$$
\begin{align*}
& {\left[\left(\rho \otimes \tau^{e}\right)(a, b) h\right](u)=\rho(a, b) h\left(\left({ }^{t} a_{v} u_{v} b_{v}\right)_{v \in \mathbf{a}}\right)}  \tag{13.19a}\\
& {\left[\left(\rho \otimes \sigma^{e}\right)(a, b) h\right](u)=\rho(a, b) h\left(\left(a_{v}^{-1} u_{v} \cdot{ }^{t} b_{v}^{-1}\right)_{v \in \mathbf{a}}\right)} \tag{13.19b}
\end{align*}
$$

for $(a, b) \in \mathfrak{K}_{0}, h \in S_{e}(T, X)$, and $u \in T$. We write these representations simply $\tau^{e}$ and $\sigma^{e}$ if $X=\mathbf{C}$ and $\rho$ is trivial.

Now we define operators $D^{e}, \bar{D}^{e}, C^{e}$, and $E^{e}$ acting on $C^{\infty}$ functions on $\mathcal{H}$ by

$$
\begin{equation*}
D^{e}=\prod_{v \in \mathbf{a}} D_{v}^{e_{v}}, \quad \bar{D}^{e}=\prod_{v \in \mathbf{a}} \bar{D}_{v}^{e_{v}}, \quad C^{e}=\prod_{v \in \mathbf{a}} C_{v}^{e_{v}}, \quad E^{e}=\prod_{v \in \mathbf{a}} E_{v}^{e_{v}} \tag{13.20}
\end{equation*}
$$

These send $C^{\infty}(\mathcal{H}, X)$ into $C^{\infty}\left(\mathcal{H}, S_{e}(T, X)\right)$. We then define an operator $D_{\rho}^{e}$ by

$$
\begin{equation*}
D_{\rho}^{e} f=\left(\rho \otimes \tau^{e}\right)(\Xi)^{-1} C^{e}[\rho(\Xi) f] \tag{13.21}
\end{equation*}
$$

Since $S_{e}(T)$ is isomorphic to $\bigotimes_{v \in \mathbf{a}} S_{e_{v}}\left(T_{v}\right)$, every irreducible subspace of $S_{e}(T)$ has multilicity 1 . Therefore, for every $\mathfrak{K}_{0}$-stable subspace $Z$ of $S_{e}(T)$ we can define a projection map $\varphi_{Z}$ of $S_{e}(T) \otimes X$ onto $Z \otimes X$. Then we identify $S_{e}(T, X)$ with $S_{e}(T) \otimes X$ by the generalization of (12.19), and define $D_{\rho}^{Z} f$ and $E^{Z}$ by

$$
\begin{equation*}
D_{\rho}^{Z} f=\varphi_{Z} D_{\rho}^{e} f, \quad E^{Z} f=\varphi_{Z} E^{e} f \tag{13.22}
\end{equation*}
$$

Denoting by $\tau_{Z}$ and $\sigma_{Z}$ the restrictions of $\tau^{e}$ and $\sigma^{e}$ to $Z$, we have obvious generalizations of (12.21), (12.22), (12.23) and (12.24a, b, c).

For $a, b \in \mathbf{Z}^{\mathbf{a}}$ let us write $a \leq b$ if $a_{v} \leq b_{v}$ for every $v \in \mathbf{a}$, and $a<b$ if $a \leq b$ and $a \neq b$. Then, as a generalization of (13.16), for $f \in \mathcal{N}^{p}(\mathcal{H}, X)$ we can put

$$
\begin{equation*}
f(z)=\sum_{a \leq p} g_{a}(z, r(z)) \tag{13.23}
\end{equation*}
$$

with a holomorphic map $g_{a}: \mathcal{H} \rightarrow S_{a}(T, X)$ for each $a$, where $g_{a}(z, u)$ means the element $g_{a}(z)$ of $S_{a}(T, X)$ evaluated at $u \in T$. Then (13.17) has the following generalization:

$$
\begin{equation*}
\left(E^{p} f\right)(u)=\left(\prod_{v \in \mathbf{a}}(-1)^{p_{v}} p_{v}!\right) g_{p}(z, u) \quad(u \in T) \tag{13.24}
\end{equation*}
$$

13.14. Lemma. Let $G$ be a Lie subgroup of $G L_{n}(\mathbf{C})$ and $L$ the Lie algebra of $G$; let $f$ be a $C^{\infty}$ map of an interval $(a, b)$ into $G$ and $\rho: G \rightarrow G L_{m}(\mathbf{C})$ be an analytic homomorphism. Then $f(t)^{-1} d f / d t \in L$ and

$$
(d / d t) \rho[f(t)]=\rho[f(t)] d \rho\left[f(t)^{-1} d f / d t\right] .
$$

Proof. For a fixed $t_{0} \in(a, b)$ and a small $\varepsilon \in \mathbf{R},>0$, define a $C^{\infty}$ map $g:(-\varepsilon, \varepsilon) \rightarrow L$ so that $g(0)=0$ and $\exp (g(h))=f\left(t_{0}\right)^{-1} f\left(t_{0}+h\right)$. Then

$$
f\left(t_{0}\right)^{-1} f^{\prime}\left(t_{0}\right)=(d / d h) \exp (g(h))_{h=0}=\lim _{h \rightarrow 0}[\exp (g(h))-1] / h=g^{\prime}(0) \in L
$$

since $g$ has values in the vector space $L$. Next,

$$
\rho\left(f\left(t_{0}+h\right)\right)=\rho\left(f\left(t_{0}\right)\right) \rho(\exp (g(h)))=\rho\left(f\left(t_{0}\right)\right) \exp (d \rho(g(h)))
$$

and hence, similarly,

$$
\frac{d}{d t}(\rho(f(t)))_{t=t_{0}}=\rho\left(f\left(t_{0}\right)\right) \frac{d}{d h} \exp (d \rho(g(h)))_{h=0}=\rho\left(f\left(t_{0}\right)\right) d \rho\left(g^{\prime}(0)\right)
$$

which combined with the equality $f\left(t_{0}\right)^{-1} f^{\prime}\left(t_{0}\right)=g^{\prime}(0)$ gives the desired result.
13.15. Proposition. (1) Let $p$ and $e$ be elements of $\mathbf{Z}^{\mathbf{a}}$ with nonnegative components such that $p_{v}=e_{v}=0$ for $v \notin \mathbf{a}^{\prime}$. If $f \in \mathcal{N}^{p}(\mathcal{H}, X)$, then $E^{e} f \in$ $\mathcal{N}^{q}\left(\mathcal{H}, S_{e}(T, X)\right)$ with $q_{v}=\operatorname{Max}\left(p_{v}-e_{v}, 0\right)$, and $D_{\rho}^{e} f \in \mathcal{N}^{p+e}\left(\mathcal{H}, S_{e}(T, X)\right)$.
(2) In particular, if $f$ is holomorphic, then

$$
\begin{equation*}
\left(D_{\rho}^{e} f\right)(z)(u)=\sum_{a<e} h_{a}(z)(u)+P_{\rho}^{e}(r(z), u) f(z) \quad(u \in T) \tag{13.25}
\end{equation*}
$$

with $h_{a} \in \mathcal{N}^{a}\left(\mathcal{H}, S_{e}(T, X)\right)$ and a map $P_{\rho}^{e}: T \times T \rightarrow \operatorname{End}(X)$ which is determined by $\rho$ and $e$ independently of $f$, and which as a function of $\left(u^{\prime}, u\right) \in T \times T$ is homogeneous of degree $e_{v}$ in $u_{v}$ and also in $u_{v}^{\prime}$ for every $v \in \mathbf{a}$.
(3) Define $p_{\rho}^{e}(r) \in S_{e}(T, \operatorname{End}(X))$ by $p_{\rho}^{e}(r)(u)=P_{\rho}^{e}(r, u)$ for $u \in T$. Then

$$
\begin{aligned}
P_{\rho}^{e+v}(r, u) x= & {\left[D_{v}\left(P_{\rho}^{e}(r, u) x\right)\right]\left(u_{v}\right)+P_{\rho}^{v}(r, u) P_{\rho}^{e}(r, u) x } \\
& +\left[d \tau^{e}\left(r_{v} \cdot{ }^{t} u_{v},{ }^{t} r_{v} u_{v}\right)\left(p_{\rho}^{e}(r) x\right)\right](u) \quad \text { for } \quad x \in X,
\end{aligned}
$$

where we view $v$ as an element of $\mathbf{Z}^{\mathbf{a}}$ such that $(v)_{v}=1$ and $(v)_{v^{\prime}}=0$ for $v^{\prime} \neq v$, and $p_{\rho}^{e}(r) x$ is the element of $S_{e}(T, X)$ such that $\left(p_{\rho}^{e}(r) x\right)(u)=p_{\rho}^{e}(r)(u) x=$ $P_{\rho}^{e}(r, u) x$.
(4) $\quad P_{\rho}^{e}\left(a u^{\prime} \cdot{ }^{t} b,{ }^{t} a^{-1} u b^{-1}\right) \rho(a, b)=\rho(a, b) P_{\rho}^{e}\left(u^{\prime}, u\right) \quad$ for every $\quad(a, b) \in \mathfrak{K}_{0}$.

Proof. The assertion concerning $E^{e} f$ is obvious in view of (13.10) and Lemma 13.3 (1). Next, take $\Xi$ as $f$ in Lemma 13.14. Since $\sum_{\nu \in N} u_{v \nu} \partial / \partial z_{v \nu}$ can be written $\sum_{k} c_{k} \partial / \partial t_{k}$ with $c_{k} \in \mathbf{C}$ and real parameters $t_{k}$, from that lemma we obtain
(13.26) $\quad \rho(\Xi)^{-1} D_{v} \rho(\Xi)\left(u^{\prime}\right)=d \rho\left(\xi^{-1}\left(D_{v} \xi\right)\left(u^{\prime}\right), \eta^{-1}\left(D_{v} \eta\right)\left(u^{\prime}\right)\right)=d \rho\left(r_{v} \cdot{ }^{t} u^{\prime},{ }^{t} r_{v} u^{\prime}\right)$
for $u^{\prime} \in T_{v}$, in view of (13.6a, b). Define $P_{\rho}^{v}: T \times T \rightarrow \operatorname{End}(X)$ by

$$
\begin{equation*}
P_{\rho}^{v}(r, u)=d \rho\left(r_{v} \cdot{ }^{t} u_{v},{ }^{t} r_{v} u_{v}\right) \tag{13.27}
\end{equation*}
$$

Then (13.26) together with (12.18) shows that

$$
\begin{equation*}
\left(D_{\rho, v} f\right)\left(u^{\prime}\right)=\left(D_{v} f\right)\left(u^{\prime}\right)+P_{\rho}^{v}\left(r, u^{\prime}\right) f \quad\left(u^{\prime} \in T_{v}\right) \tag{13.28}
\end{equation*}
$$

If $f$ is a polynomial in $r$ of degree $\leq p$ with holomorphic coefficients, then (13.6c) shows that $D_{v} f$ is a polynomial in $r$ of degree $\leq p+v$. The same is true for $D_{\rho, v} f$, since $P_{\rho}^{v}(r, u)$ is linear in $r$. By Proposition 12.10 (1) and induction, we see that $D_{\rho}^{e} f \in \mathcal{N}^{p+e}\left(\mathcal{H}, S_{e}(T, X)\right)$. We prove (2) by induction on $\sum_{v \in \mathbf{a}} e_{v}$. Our assertion is trivial if $e=0$. Now fix one $v$, and assume (13.25) for some $e$; put $h=\sum_{a<e} h_{a}$. and define $p_{\rho}^{e}(r)$ as in (3). By Proposition 12.10, $D_{\rho}^{e+v} f=D_{\rho \otimes \tau^{e}, v}\left(D_{\rho}^{e} f\right)$, and so

$$
\left(D_{\rho}^{e+v} f\right)\left(u^{\prime}\right)=\left(D_{\rho \otimes \tau^{e}, v} h\right)\left(u^{\prime}\right)+\left[D_{\rho \otimes \tau^{e}, v}\left(p_{\rho}^{e}(r) f\right)\right]\left(u^{\prime}\right)
$$

Here we view $p_{\rho}^{e}(r) f$ as an element of $C^{\infty}\left(\mathcal{H}, S_{e}(T, X)\right)$ by the rule $\left(p_{\rho}^{e}(r) f\right)(u)=$ $p_{\rho}^{e}(r)(u) f$. Now $D_{\rho \otimes \tau^{e}, v} h \in \sum_{b<e+v} \mathcal{N}^{b}\left(\mathcal{H}, S_{1}\left(T_{v}, S_{e}(T, X)\right)\right)$ by (1). By (13.28) we have

$$
\begin{aligned}
{\left[D_{\rho \otimes \tau^{e}, v}\left(p_{\rho}^{e}(r) f\right)\right]\left(u^{\prime}\right) } & =D_{v}\left(p_{\rho}^{e}(r) f\right)\left(u^{\prime}\right)+P_{\rho \otimes \tau^{e}}^{v}\left(r, u^{\prime}\right)\left(p_{\rho}^{e}(r) f\right) \\
& =p_{\rho}^{e}(r)\left(D_{v} f\right)\left(u^{\prime}\right)+Q\left(r, u^{\prime}\right) f
\end{aligned}
$$

for $u^{\prime} \in T_{v}$ with an element $Q\left(r, u^{\prime}\right)$ of $S_{1}\left(X, S_{e}(T, X)\right)$ given by $Q\left(r, u^{\prime}\right) x=$ $D_{v}\left(p_{\rho}^{e}(r) x\right)\left(u^{\prime}\right)+P_{\rho \otimes \tau^{e}}^{v}\left(r, u^{\prime}\right)\left(p_{\rho}^{e}(r) x\right)$ for $x \in X$. Since $p_{\rho}^{e}(r)$ is homogeneous of degree $e$ in $r$, from (13.6c) we see that $D_{v}\left(p_{\rho}^{e}(r)\right)$ is homogeneous of degree $e+v$ in $r$. Also, $P_{\rho \otimes \tau^{e}}^{v}\left(r, u^{\prime}\right)$ is bilinear in ( $\left.r, u^{\prime}\right)$. Thius we obtain (13.25) for $D_{\rho}^{e+v} f$ with $P_{\rho}^{e+v}(r, u) x=\left(Q\left(r, u_{v}\right) x\right)(u)$, that is,

$$
P_{\rho}^{e+v}(r, u) x=\left[D_{v}\left(P_{\rho}^{e}(r, u) x\right)\right]\left(u_{v}\right)+\left[P_{\rho \otimes \tau^{e}}^{v}(r, u)\left(p_{\rho}^{e}(r) x\right)\right](u) \quad(u \in T, x \in X) .
$$

This proves (2). To prove (3), take $h \in S_{e}(T, X)$; then

$$
\begin{aligned}
& {\left[d\left(\rho \otimes \tau^{e}\right)(A, B) h\right](u)} \\
& \quad=(d / d t)_{t=0}\left\{\rho(\exp (t A), \exp (t B)) h\left({ }^{t} \exp (t A) \cdot u \cdot \exp (t B)\right)\right\} \\
& \quad=d \rho(A, B) \cdot h(u)+\left[d \tau^{e}(A, B) h\right](u)
\end{aligned}
$$

Take $h=p_{\rho}^{e}(r) x$. Then

$$
\begin{aligned}
& {\left[P_{\rho \otimes \tau^{e}}^{v}(r, u)\left(p_{\rho}^{e}(r) x\right)\right](u)} \\
& \quad=P_{\rho}^{v}(r, u)\left(p_{\rho}^{e}(r) x\right)(u)+\left[d \tau^{e}\left(r_{v} \cdot{ }^{t} u_{v},{ }^{t} r_{v} u_{v}\right)\left(p_{\rho}^{e}(r) x\right)\right](u)
\end{aligned}
$$

This proves (3). Finally, combining (13.24) and (13.25), we obtain $\left(E^{e} D_{\rho}^{e} f\right)(u, w)$ $=c P_{\rho}^{e}(w, u) f(z)$ with $c=\prod_{v \in \mathbf{a}}(-1)^{e_{v}} e_{v}!$. Now replace $f$ by $f \|_{\rho} \alpha$ with $\alpha \in G_{\mathbf{a}}$. By Proposition 12.10 (2), ( $\left.E^{e} D_{\rho}^{e} f\right) \|_{\rho \otimes \tau^{e} \otimes \sigma^{e}} \alpha=E^{e} D_{\rho}^{e}\left(f \|_{\rho} \alpha\right)$, and hence

$$
P_{\rho}^{e}(w, u)\left(f \|_{\rho} \alpha\right)=\rho\left(\lambda_{\alpha}, \mu_{\alpha}\right)^{-1} P_{\rho}^{e}\left(\lambda_{\alpha} w \cdot{ }^{t} \mu_{\alpha},{ }^{t} \lambda_{\alpha}^{-1} u \mu_{a}^{-1}\right)(f \circ \alpha)
$$

for every holomorphic map $f: \mathcal{H} \rightarrow X$. Therefore

$$
P_{\rho}^{e}(w, u) \rho\left(\lambda_{\alpha}, \mu_{\alpha}\right)^{-1}=\rho\left(\lambda_{\alpha}, \mu_{\alpha}\right)^{-1} P_{\rho}^{e}\left(\lambda_{\alpha} w \cdot{ }^{t} \mu_{\alpha},{ }^{t} \lambda_{\alpha}^{-1} u \mu_{a}^{-1}\right),
$$

from which we obtain (4).
13.16. Lemma. Given $\{\rho, X\}$ as before and $f \in C^{\infty}(\mathcal{H}, X)$, for $u, u^{\prime} \in T_{v}$ with a fixed $v \in \mathbf{a}^{\prime}$, we have

$$
\begin{gathered}
{\left[\left(D_{\rho \otimes \sigma_{v}, v} E_{v}-E_{v} D_{\rho, v}\right) f\right]\left(u, u^{\prime}\right)=P_{\rho}^{v}\left(u^{\prime}, u\right) f} \\
\left(L_{\rho, v}-M_{\rho, v}\right) f=B_{\rho, v} f
\end{gathered}
$$

with a constant element $B_{\rho, v}$ of $\operatorname{End}(X)$ depending only on $\rho$ and $v$, where $L_{\rho, v}=-\theta D_{\rho \otimes \sigma_{v}, v} E_{v}$ and $M_{\rho, v}=-\theta E_{v} D_{\rho, v}$ with $\theta$ of (12.33), and $u$ (resp. $u^{\prime}$ ) corresponds to the operator $D_{v}$ (resp. $E_{v}$ ).

Proof. By (12.12b) and (12.18) we have

$$
\begin{aligned}
\left(D_{\rho \otimes \sigma_{v}, v} E_{v} f\right)\left(u, u^{\prime}\right) & =\left\{\left(\rho \otimes \sigma_{v}\right)(\Xi)^{-1} D_{v}\left[\rho(\Xi) \bar{D}_{v} f\left(u^{\prime}\right)\right]\right\}(u) \\
& =\rho(\Xi)^{-1} D_{v}(\rho(\Xi))(u)\left(\bar{D}_{v} f\right)\left(\xi u^{\prime} \cdot t \eta\right)+\left(D_{v} \bar{D}_{v} f\right)\left(\xi u^{\prime} \cdot t \eta, u\right)
\end{aligned}
$$

By (13.26) the first term of the last line equals $P_{\rho}^{v}(r, u)\left(E_{v} f\right)\left(u^{\prime}\right)$. On the other hand, by (13.28),

$$
\left(E_{v} D_{\rho, v} f\right)\left(u, u^{\prime}\right)=E_{v}\left[\left(D_{v} f\right)(u)\right]\left(u^{\prime}\right)+E_{v}\left[P_{\rho}^{v}(r, u) f\right]\left(u^{\prime}\right) .
$$

By (12.12b) the first term on the right-hand side is $\left(\bar{D}_{v} D_{v} f\right)\left(u, \xi u^{\prime} \cdot{ }^{t} \eta\right)$, and by (13.10) the second term equals $-P_{\rho}^{v}\left(u^{\prime}, u\right) f+P_{\rho}^{v}(r, u)\left(E_{v} f\right)\left(u^{\prime}\right)$. Since $\left(\bar{D}_{v} D_{v} f\right)$ $\left(u, u^{\prime}\right)=\left(D_{v} \bar{D}_{v} f\right)\left(u^{\prime}, u\right)$, we obtain the first formula, which immediately implies the second formula with $B_{\rho, v}=-\theta P_{\rho}^{v}$.
13.17. Lemma. Let $\{\rho, X\}$ and $\{\sigma, Y\}$ be representations of $\mathfrak{K}_{0}$, and let $f \in C^{\infty}(\mathcal{H}, X)$ and $g \in C^{\infty}(\mathcal{H}, Y)$. Then for $u \in T$ we have

$$
\begin{equation*}
D_{\rho \otimes \sigma}^{e}(f \otimes g)(u)=\sum_{a+b=e}\binom{e}{a}\left(D_{\rho}^{a} f\right)(u) \otimes\left(D_{\sigma}^{b} g\right)(u) \tag{13.29}
\end{equation*}
$$

where the sum is extended over all $a, b \in \mathbf{Z}^{\mathbf{a}}$ with nonnegative components such that $a+b=e$, and $\binom{e}{a}=\prod_{v \in \mathbf{a}} \frac{e_{v}!}{a_{v}!\left(e_{v}-a_{v}\right)!}$. Similarly if $X=Y=\operatorname{End}(W)$, $\rho(a, b) x=\rho_{0}(a, b) x$ for $x \in \operatorname{End}(W)$ with a representation $\left\{\rho_{0}, W\right\}$, and $\sigma$ is trivial, then the above formula holds with $\otimes$ replaced by multiplication in $\operatorname{End}(W)$.

Proof. Since $C_{v}$ is given by (13.10), for $f_{1} \in C^{\infty}(\mathcal{H}, X)$ and $g_{1} \in C^{\infty}(\mathcal{H}, Y)$ we have $C_{v}\left(f_{1} \otimes g_{1}\right)(u)=\left(C_{v} f_{1}\right)(u) \otimes g_{1}+f_{1} \otimes\left(C_{v} g_{1}\right)(u)$ for $u \in T_{v}$, and

$$
\begin{aligned}
C_{v^{\prime}} C_{v}\left(f_{1} \otimes g_{1}\right)\left(u, u^{\prime}\right)= & \left(C_{v^{\prime}} C_{v} f_{1}\right)\left(u, u^{\prime}\right) \otimes g_{1}+\left(C_{v} f_{1}\right)(u) \otimes\left(C_{v^{\prime}} g_{1}\right)\left(u^{\prime}\right) \\
& +\left(C_{v^{\prime}} f_{1}\right)\left(u^{\prime}\right) \otimes\left(C_{v} g_{1}\right)(u)+f_{1} \otimes\left(C_{v^{\prime}} C_{v} g_{1}\right)\left(u, u^{\prime}\right)
\end{aligned}
$$

for $u \in T_{v}$ and $u^{\prime} \in T_{v^{\prime}}$. Applying the $C_{v}$ successively, for $u \in T$ we find that

$$
C^{e}\left(f_{1} \otimes g_{1}\right)(u)=\sum_{a+b=e}\binom{e}{a}\left(C^{a} f_{1}\right)(u) \otimes\left(C^{b} g_{1}\right)(u) .
$$

By (13.21), $D_{\rho \otimes \sigma}^{e}(f \otimes g)(u)=(\rho \otimes \sigma)(\Xi)^{-1} C^{e}(\rho(\Xi) f \otimes \sigma(\Xi) g)\left({ }^{t} \xi^{-1} u \eta^{-1}\right)$, and therefore taking $f_{1}=\rho(\Xi) f$ and $g_{1}=\sigma(\Xi) g$, we obtain (13.29). The case $X=$ $Y=\operatorname{End}(W)$ can be proved in the same manner.

## 14. Arithmeticity of nearly holomorphic functions

14.0. Before discussing our problems in the general case, let us first illustrate the main ideas by sketching the proof of the following statement in which we take $G=S L_{2}(\mathbf{Q}):$

Let $f \in \mathcal{M}_{k}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ and $g \in \mathcal{A}_{k+2 p}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ with positive integers $k$ and $p$. Let $K$ be an imaginary quadratic field embedded in $\mathbf{C}$, and let $\tau \in H \cap K$; suppose that $g$ has neither pole nor zero at $\tau$. Then $\left(\pi^{-p} \delta_{k}^{p} f / g\right)(\tau)$ belongs to $K_{\mathrm{ab}}$, where $\delta_{k}^{p}$ is the operator of (12.39), which equals $D_{\rho}^{p}$ with $\rho(x)=x^{k}$ as explained in $\S 12.17$.

This is one of the easiest cases of (0.6) in the introduction, formulated with $K_{\text {ab }}$ instead of $\overline{\mathbf{Q}}$. Our proof is by induction on $p$. We may assume that $f(\tau) \neq 0$. Indeed, if $f(\tau)=0$, then we take $f_{1} \in \mathcal{M}_{k}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ so that $f_{1}(\tau) \neq 0$, and put $f_{2}=$ $f+f_{1}$. Then our assertion on $\left(\pi^{-p} \delta_{k}^{p} f / g\right)(\tau)$ follows from that on $\left(\pi^{-p} \delta_{k}^{p} f_{\nu} / g\right)(\tau)$ for $\nu=1,2$. We first observe that the group $\left\{\alpha \in S L_{2}(\mathbf{Q}) \mid \alpha(\tau)=\tau\right\}$ contains an element, say $\alpha$, of infinite order. Assuming $f(\tau) \neq 0$, put $h=\left(f \|_{k} \alpha\right) / f$. Then $h$ is a $\mathbf{Q}_{\mathrm{ab}}$-rational modular function, so that $h(\tau) \in K_{\mathrm{ab}}$ by classical theory of complex multiplication, which is actually a special case of Theorem 9.6. Now we easily see that $\delta_{k}(r s)=r^{\prime} s+r \delta_{k} s$ for any meromorphic functions $r$ and $s$, where $r^{\prime}=d r / d z$. Applying this principle to the equality $f \|_{k} \alpha=h f$, we obtain $\left(\delta_{k} f\right) \|_{k+2} \alpha=\delta_{k}\left(f \|_{k} \alpha\right)=\delta_{k}(h f)=h^{\prime} f+h \delta_{k} f$. Put $\lambda=j_{\alpha}(\tau)$. Then $\lambda$ is an eigenvalue of $\alpha$, so that $\lambda$ is an element of $K$ that is not a root of unity. We have also $\left(f \|_{k} \alpha\right)(\tau)=\lambda^{-k} f(\tau), h(\tau)=\lambda^{-k}$, and $\left(\left(\delta_{k} f\right) \|_{k+2} \alpha\right)(\tau)=$ $\lambda^{-k-2}\left(\delta_{k} f\right)(\tau)$. Therefore $\left(\lambda^{-2}-1\right)\left(\delta_{k} f\right)(\tau)=\lambda^{k} h^{\prime}(\tau) f(\tau)$. Now $\pi^{-1} h^{\prime} \in \mathcal{A}_{2}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, and hence $\left(\pi^{-1} h^{\prime} f / g\right) \in \mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, if $g \in \mathcal{A}_{k+2}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ and $g(\tau) \neq 0$. We have thus $\left(\pi^{-1}\left(\delta_{k} f\right) / g\right)(\tau)=\lambda^{k}\left(\lambda^{-2}-1\right)^{-1}\left(\pi^{-1} h^{\prime} f / g\right)(\tau) \in K_{\mathrm{ab}}$, which proves the case $p=1$.

If $p>1$, we first observe that $\delta_{k}^{p}(r s)=\sum_{a=0}^{p} c_{a}\left(\delta_{0}^{a} r\right)\left(\delta_{k}^{p-a} s\right)$ with $c_{a}=$ $p!/[a!(p-a)!]$, and so

$$
\left(\delta_{k}^{p} f\right) \|_{k+2 p} \alpha=\delta_{k}^{p}\left(f \|_{k} \alpha\right)=\delta_{k}^{p}(h f)=h \delta_{k}^{p} f+\sum_{a=1}^{p} c_{a}\left(\delta_{0}^{a} h\right)\left(\delta_{k}^{p-a} f\right)
$$

Evaluating this equality at $\tau$, we obtain

$$
\left(\delta_{k}^{p} f\right)(\tau)\left(\lambda^{-2 p}-1\right) \lambda^{-k}=\sum_{a=1}^{p} c_{a}\left(\delta_{0}^{a} h\right)(\tau) \cdot\left(\delta_{k}^{p-a} f\right)(\tau) .
$$

Take $q_{a} \in \mathcal{A}_{2 a}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ so that $q_{a}(\tau) \neq 0$. Then

$$
\pi^{-p}\left(\delta_{k}^{p} f / g\right)(\tau)=\left(\lambda^{-2 p}-1\right)^{-1} \lambda^{k} \sum_{a=1}^{p} c_{a}\left(\pi^{-a} \delta_{0}^{a} h / q_{a}\right)(\tau) \cdot\left(\pi^{a-p} q_{a} \delta_{k}^{p-a} f / g\right)(\tau)
$$

Since $g / q_{a} \in \mathcal{A}_{k-a+2 p}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, our induction shows that $\left(\pi^{a-p} q_{a} \delta_{k}^{p-a} f / g\right)(\tau) \in K_{\mathrm{ab}}$. As for the other factor, we have $\pi^{-a} \delta_{0}^{a} h=\pi^{1-a} \delta_{2}^{a-1}\left(\pi^{-1} h^{\prime}\right)$, and $\pi^{-1} h^{\prime} \in \mathcal{A}_{2}\left(\mathbf{Q}_{\mathrm{ab}}\right)$. Therefore our induction is applicable to that factor. We can thus conclude that $\pi^{-p}\left(\delta_{k}^{p} f / g\right)(\tau) \in K_{\text {ab }}$.

We shall prove in Theorem 14.7 below a generalization of the above statement, essentially by the same idea, and then in Theorem 14.9 a similar but much stronger result. Since we have to deal with vector-valued or matrix-valued functions, our treatment will become more involved than the above proof. The generalization of the fact $\pi^{-1} h^{\prime} \in \mathcal{A}_{2}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, given in Proposition 14.5, is highly nontrivial in Case UB.
14.1. Thus our principal interest in this section is in the nature of the value of $D_{\rho}^{Z} f$ at a CM-point of $\mathcal{H}$. We start with some results without arithmeticity.

Let $\{\rho, X\}$ be a representation of $\mathfrak{K}_{0}$. Given $h \in S_{e}\left(T, S_{e}(T, X)\right)$ and $(u, w) \in$ $T \times T$, we define $h_{u}, h^{w} \in S_{e}(T, X)$ by $h_{u}(w)=h^{w}(u)=h(u, w)$. Then we obtain $p, q \in S_{e}\left(T, M l_{e}(T, X)\right)$ defined by $p(u)=\left(h_{u}\right)_{*}$ and $q(w)=\left(h^{w}\right)_{*}$. Moreover, $p_{*}$ and $q_{*}$ are meaningful as elements of $M l_{e}\left(T, M l_{e}(T, X)\right)$. We can easily verify that $p_{*}=q_{*}$; we then denote this same element of $M l_{e}\left(T, M l_{e}(T, X)\right)$ by $h_{* *}$, and define $\theta_{X}^{e}: S_{e}\left(T, S_{e}(T, X)\right) \rightarrow X$ by

$$
\begin{equation*}
\theta_{X}^{e} h=\sum h_{* *}\left(a_{\nu_{1}}, a_{\nu_{2}}, \ldots ; b_{\nu_{1}}, b_{\nu_{2}}, \ldots\right) \tag{14.1}
\end{equation*}
$$

Here we write $T^{e}=\prod_{v \in \mathbf{a}} T_{v}^{e_{v}}=\prod_{i=1}^{|e|} T_{i},|e|=\sum_{v \in \mathbf{a}} e_{v}$, where each $T_{i}$ is identified with some $T_{v} ;\left\{a_{\nu_{i}}\right\}$ and $\left\{b_{\nu_{i}}\right\}$ are dual bases of $T_{i}$; the summation is over all the elements $\left\{\left(a_{\nu_{i}}, b_{\nu_{i}}\right)\right\}_{i=1}^{|e|}$. Then we can easily verify that

$$
\begin{equation*}
\theta_{X}^{e} \circ\left(\rho \otimes \tau^{e} \otimes \sigma^{e}\right)(\alpha)=\theta_{X}^{e} \circ\left(\rho \otimes \sigma^{e} \otimes \tau^{e}\right)(\alpha)=\rho(\alpha) \circ \theta_{X}^{e} \quad \text { for every } \quad \alpha \in \mathfrak{K}_{0} . \tag{14.2}
\end{equation*}
$$

Given $g \in C^{\infty}\left(\mathcal{H}, S_{e}(T, X)\right)$, we can view $D_{\rho \otimes \sigma^{e}}^{e} g$ as an element of $C^{\infty}(\mathcal{H}$, $S_{e}\left(T, S_{e}(T, X)\right)$. Thus $\theta_{X}^{e} D_{\rho \otimes \sigma^{e}}^{e} g$ is meaningful as an element of $C^{\infty}(\mathcal{H}, X)$.
14.2. Proposition. Let $\{\rho, X\}$ and $\left\{\rho_{0}, X\right\}$ be representations of $\mathfrak{K}_{0}$ such that $\rho(a, b)=\operatorname{det}(b)^{k} \rho_{0}(a, b)$ for $(a, b) \in \mathfrak{K}_{0}$ with $k \in \mathbf{Z}^{\mathbf{a}}$. Let $f$ be an element of $\mathcal{N}^{p}(\mathcal{H}, X)$ such that $f \|_{\rho} \gamma=f$ for every $\gamma$ in a subgroup $\Gamma$ of $G_{\mathbf{a}}$. Then there exists an integer $N\left(\rho_{0}, p\right)$ that depends only on $\rho_{0}$ and $p$ with the following property: if $k_{v}>N\left(\rho_{0}, p\right)$ for every $v \in \mathbf{a}^{\prime}$, then

$$
f=\sum_{e \leq p} \theta_{X}^{e} D_{\rho \otimes \sigma^{e}}^{e} g_{e}
$$

with holomorphic maps $g_{e}: \mathcal{H} \rightarrow S_{e}(T, X)$ such that $g_{e} \|_{\rho \otimes \sigma^{e}} \gamma=g_{e}$ for every $\gamma \in \Gamma$, where $e$ runs over the elements of $\mathbf{Z}^{\mathbf{a}}$ such that $0 \leq e_{v} \leq p_{v}$ for every $v \in \mathbf{a}$. Here $\mathbf{a}^{\prime}=\left\{v \in \mathbf{a} \mid G_{v}\right.$ is not compact $\}$ as already defined in (12.5).

Proof. We prove this by induction on $|p|$. Our assertion is trivial if $p=0$. Fixing $e$, put $\pi_{0}=\rho_{0} \otimes \sigma^{e}$ and $\pi(a, b)=\operatorname{det}(b)^{k} \pi_{0}(a, b)$, and $Y_{e}=S_{e}(T, X)$. Then $d \pi(x, y)=d \pi_{0}(x, y)+\sum_{v \in \mathbf{a}} k_{v} \cdot \operatorname{tr}\left(y_{v}\right) 1_{Y}$, where $1_{Y}$ is the identity element of $\operatorname{End}\left(Y_{e}\right)$. Define $P_{\rho}^{e}$ and $P_{\rho}^{v}$ by (13.25) and (13.27) (with $Y_{e}$ as $X$ there). Then $P_{\pi}^{v}(r, u)=P_{\pi_{0}}^{v}(r, u)+\sum_{v} k_{v} \cdot \operatorname{tr}\left({ }^{t} r_{v} u_{v}\right) 1_{Y}$. Now we have

$$
\begin{align*}
& P_{\pi}^{e}(w, u)=\sum_{\nu \leq e} k^{\nu} q_{\nu}(w, u) \text { with } q_{\nu} \in S_{e}\left(T, S_{e}\left(T, \operatorname{End}\left(Y_{e}\right)\right)\right)  \tag{14.3}\\
& \text { depending only on } \rho_{0} \text { and } e \text {; moreover } q_{e}(w, u)=\operatorname{tr}\left({ }^{t} w u\right)^{e} 1_{Y} .
\end{align*}
$$

We prove this by induction on $|e|=\sum_{v \in \mathbf{a}} e_{v}$. If $|e|=1$, we have $P_{\pi}^{e}(w, u)=$ $P_{\pi}^{v}(w, u)$ with some $v$, and $q_{e}(w, u)=\operatorname{tr}\left({ }^{t} w_{v} u_{v}\right) 1_{Y}$; thus (14.3) is true if $|e|=1$. Assuming (14.3) for some $e$, define $p_{\pi}^{e}(w), q_{\nu}(w) \in S_{e}\left(T, \operatorname{End}\left(Y_{e}\right)\right)$ by $p_{\pi}^{e}(w)(u)=$ $P_{\pi}^{e}(w, u)$ and $q_{\nu}(w)(u)=q_{\nu}(w, u)$ for $w, u \in T$. Then, by Proposition 13.15 (3), for $y \in Y_{e}$ we have

$$
\begin{aligned}
P_{\pi}^{e+v}(r, u) y & =\left[D_{v}\left(\sum_{\nu \leq e} k^{\nu} q_{\nu}(r, u) y\right)\right]\left(u_{v}\right)+P_{\pi}^{v}(r, u) \sum_{\nu \leq e} k^{\nu} q_{\nu}(r, u) y \\
& +\left[d \tau^{e}\left(r_{v} \cdot{ }^{t} u_{v},{ }^{t} r_{v} u_{v}\right)\left(\sum_{\nu \leq e} k^{\nu} q_{\nu}(r) y\right)\right](u) .
\end{aligned}
$$

Thus we easily see that $P_{\pi}^{e+v}(r, u)=\sum_{\nu \leq e+v} k^{\nu} s_{\nu}(r, u)$ with $s_{\nu} \in S_{e}\left(T, S_{e}(T\right.$, $\left.\operatorname{End}\left(Y_{e}\right)\right)$ ) depending only on $\rho_{0}$ and $e$. Moreover $s_{e+v}(r, u)=\operatorname{tr}\left({ }^{t} r_{v} u_{v}\right) q_{e}(r, u)$. Thus we obtain (14.3) for $P_{\pi}^{e+v}$.

Next, define $A_{e}, B_{e . \nu} \in \operatorname{End}\left(Y_{e}\right)$ by

$$
\left(A_{\epsilon} \varphi\right)(w)=\theta_{X}^{e}\left(\left[P_{\pi}^{e}(w, *) \varphi\right](*)\right) \text { and }\left(B_{e . \nu} \varphi\right)(w)=\theta_{X}^{e}\left(\left[q_{\nu}(w, *) \varphi\right](*)\right)
$$

for $\varphi \in Y_{e}=S_{e}(T, X)$ and $w \in T$, where $[q(w, *) \varphi](*)$ is an element of $S_{e}\left(T, Y_{e}\right)=$ $S_{e}\left(T, S_{e}(T, X)\right)$ whose value at $\left(u_{1}, u_{2}\right) \in T \times T$ is $\left[q\left(w, u_{1}\right) \varphi\right]\left(u_{2}\right)$. We have $q_{e}(w, u)=\operatorname{tr}\left({ }^{t} w u\right)^{e} 1_{Y}$, so that (12.10) implies that $\left(B_{e . \epsilon} \varphi\right)(w)=\varphi(w)$, that is, $B_{e . e}=1_{Y}$. Thus $A_{\epsilon}=k^{e} 1_{Y}+\sum_{\nu<e} k^{\nu} B_{e . \nu}$. Observe that $A_{e}$ is invertible if $k_{v}>$ $M\left(\rho_{0}, e\right)$ for every $v \in \mathbf{a}^{\prime}$ with an integer $M\left(\rho_{0}, e\right)$ that depends only on $\rho_{0}$ and $e$. Observe also that if a map $g: \mathcal{H} \rightarrow Y_{e}$ is holomorphic, then $D_{\pi}^{e} g=\sum_{a \leq e} h_{a}(r)$ with holomorphic maps $h_{a}: \mathcal{H} \rightarrow S_{a}\left(T, S_{e}\left(T, Y_{e}\right)\right)$, and $\theta_{X}^{e} D_{\pi}^{e} g=\sum_{a \leq \epsilon} \theta_{X}^{e} h_{a}(r)$. By
(13.25), $h_{e}(r)=P_{\pi}^{e}(r, *) g$, and so $\theta_{X}^{e} h_{e}(r)=\left(A_{e} g\right)(r)$. Now we consider a given element $f$ of $\mathcal{N}_{\rho}^{p}$ of our proposition. Then $f=\sum_{a \leq p} \ell_{a}(r)$ with holomorphic maps $\ell_{a}: \mathcal{H} \rightarrow S_{a}(T, X)$ as noted in (13.23). By (13.24), $\ell_{p}=c E^{p} f$ with $c \in \mathbf{Q}^{\times}$, and so $\ell_{p} \|_{\rho \otimes \sigma^{p}} \gamma=\ell_{p}$ for every $\gamma \in \Gamma$ by Proposition 12.10 (2). Define an integer $N\left(\rho_{0}, p\right)$ by $N\left(\rho_{0}, p\right)=\operatorname{Max}_{e \leq p}\left\{M\left(\rho_{0}, e\right)\right\}$. Suppose $k_{v}>N\left(\rho_{0}, p\right)$ for every $v \in \mathbf{a}^{\prime}$ and put $g=A_{p}^{-1} \ell_{p}$. (So we take $p$ as the above $e$ and consider $P_{\pi}^{p}$ with $\pi=\rho \otimes \sigma^{p}$.) From Proposition 13.15 (4) and (14.2) we see that $A_{p}$ commutes with $\left(\rho \otimes \sigma^{p}\right)\left(\mathfrak{K}_{0}\right)$, and hence $g$ is a holomorphic map of $\mathcal{H}$ into $S_{p}(T, X)$ such that to $g \|_{\rho \otimes \sigma^{p}} \gamma=g$ for every $\gamma \in \Gamma$. Now the $p$-th degree part of $\theta_{X}^{p} D_{\rho \otimes \sigma^{p}}^{p} g$ is $\left(A_{p} g\right)(r)=\ell_{p}(r)$, so that $f-\theta_{X}^{p} D_{\rho \otimes \sigma^{p}}^{p} g \in \sum_{a<p} \mathcal{N}^{a}(\mathcal{H}, X)$. Applying our induction to the last difference, we complete the proof.
14.3. Lemma. $\mathcal{N}_{\rho}^{p}(\Gamma)$ is finite-dimensional over $\mathbf{C}$ for every congruence subgroup $\Gamma$ of $G$.

Proof. We can find a positive integer $\kappa$ and a nonzero element $h$ of $\mathcal{M}_{\kappa \mathbf{a}}\left(\Gamma^{\prime}\right)$ with a congruence subgroup $\Gamma^{\prime}$ contained in $\Gamma$. Such an $h$ can be obtained by $h=\prod_{v \in \mathbf{a}} \operatorname{det}\left(Q_{v}\right)$ with $Q_{v}$ of Proposition 11.14. Take a positive integer $m$ so that $m \kappa>N(\rho, p)$ with $N(\rho, p)$ as in Proposition 14.2. Given $f \in \mathcal{N}_{\rho}^{p}(\Gamma)$, we have $h^{m} f \in \mathcal{N}_{\rho^{\prime}}^{p}\left(\Gamma^{\prime}\right)$, where $\rho^{\prime}(a, b)=\rho(a, b) \operatorname{det}(b)^{m \kappa \mathbf{a}}$. By Proposition 14.2, $h^{m} f=\sum_{e \leq p} \theta_{X}^{e} D_{\rho^{\prime} \otimes \sigma^{e}}^{e} g_{e}$ with $g_{e} \in \mathcal{M}_{\rho^{\prime} \otimes \sigma^{e}}\left(\Gamma^{\prime}\right)$. Since $\mathcal{M}_{\rho^{\prime} \otimes \sigma^{e}}\left(\Gamma^{\prime}\right)$ is finitedimensional, we see that $h^{m} f$ belongs to a finite-dimensional vector space over $\mathbf{C}$, which implies our lemma.
14.4. To consider arithmeticity, we go back to the setting of Sections 5 and 11, and note that $T_{v}$ has a natural $\overline{\mathbf{Q}}$-structure. In each case the structure can be obtained by taking a natural coordinate system of $T_{v}$ determined by the matrix entries. In Case UB, as already noted in the proof of Lemma 4.13, the action of $\widetilde{G}_{+}$on $\mathcal{H}$ is $\overline{\mathbf{Q}}$-rational.

Let us now take a $\overline{\mathbf{Q}}$-rational representation $\{\omega, X\}$ of $\mathfrak{K}$, where

$$
\begin{equation*}
\mathfrak{K}=\prod_{v \in \mathbf{b}} G L_{n_{v}}(\mathbf{C}) . \tag{14.4}
\end{equation*}
$$

Since $\mathfrak{K}_{0} \subset \mathfrak{K}$, we can speak of the restriction of $\omega$ to $\mathfrak{K}_{0}$. Taking this restriction to be $\rho$, we can define various objects with respect to $\rho$. To define arithmeticity, however, we have to consider them relative to $\omega$. First of all, we can express an element of $\mathfrak{K}$ in the form ( $a, b$ ) with $a \in \prod_{v \in \mathbf{a}} G L_{m_{v}}(\mathbf{C})$ and $b \in \prod_{v \in \mathbf{a}} G L_{n_{v}}(\mathbf{C})$; for Type C we have $\mathfrak{K}=\mathfrak{K}_{0}$ and $\omega=\rho$, and we take $a=b$ and $m_{v}=n_{v}=n$; for Type A if $\alpha=\left(\alpha_{v}\right)_{v \in \mathbf{b}} \in \mathfrak{K}$, then $a=\left(\alpha_{v \rho}\right)_{v \in \mathbf{a}}$ and $b=\left(\alpha_{v}\right)_{v \in \mathbf{a}}$; see $\S 5.1$ for our convention. We can then define $\omega \otimes \tau^{e}$ and $\omega \otimes \sigma^{e}$ by (13.19a, b) with $\omega$ in place of $\rho$. We denote the symbols $D_{\rho}^{e}, D_{\rho}^{Z}, \mathcal{N}_{\rho}^{p}$, etc. by $D_{\omega}^{e}, D_{\omega}^{Z}, \mathcal{N}_{\omega}^{p}$, etc. Then (12.21) is true for $\alpha$ in $G=U(\mathcal{T})$ or $G=U\left(\eta_{n}\right)$ in the unitary case with $Z \subset S_{e}(T)$ and $\omega$ in place of $\rho$. This is because (12.1a, b) are true for such an $\alpha$, and so the proof of Proposition 12.10 is valid. If $\omega(x)=\operatorname{det}(x)^{k}$ with $k \in \mathbf{Z}^{\mathbf{b}}$, then we replace the subscript $\omega$ by $k$.

Let $f \in C^{\infty}(\mathcal{H}, X)$ and $w \in \mathcal{H}_{\mathrm{CM}}$. Then we say that $f$ is $\omega$-arithmetic (or simply arithmetic, if $\omega$ is clear from the context) at $w$ if $\mathfrak{P}_{\omega}(w)^{-1} f(w)$ is $\overline{\mathbf{Q}}$ rational. Then we denote by $\mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$ (resp. $\mathcal{N}_{\omega}^{p}(\Gamma, \overline{\mathbf{Q}})$ ) the set of the elements of $\mathcal{N}_{\omega}^{p}$ (resp. $\left.\mathcal{N}_{\omega}^{p}(\Gamma)\right)$ that are $\omega$-arithmetic at every point of $\mathcal{H}_{\mathrm{CM}}$. If $\{\omega, X\}$ is clear from the context, we simply call an element of $\mathcal{N}_{\omega}^{p}$ arithmetic if it belongs to $\mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$. From Proposition 11.5 (2) we see that if $f \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$ and $\alpha \in \widetilde{G}_{+}$, then
$f \|_{\omega} \alpha \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$. For Type C we have $\omega=\rho$, and so $\mathcal{N}_{\rho}^{p}(\overline{\mathbf{Q}})$ is meaningful. For Type A, however, we cannot speak of $\mathcal{N}_{\rho}^{p}(\overline{\mathbf{Q}})$ for the reasons explained in Theorem 11.17.
14.5. Proposition. For each $v \in \mathbf{a}^{\prime}$ define a representation $\left\{\tau_{v}, S_{1}\left(T_{v}\right)\right\}$ of $\mathfrak{K}$ by $\left[\tau_{v}(a, b) h\right](u)=h\left({ }^{t} a_{v} u b_{v}\right)$ for $h \in S_{1}\left(T_{v}\right)$ and $u \in T_{v}$. (This is a special case of (12.7a).) Then $\pi^{-1} D_{v} f \in \mathcal{A}_{\tau_{v}}(\overline{\mathbf{Q}})$ for every $f \in \mathcal{A}_{0}(\overline{\mathbf{Q}})$.

Proof. By (12.14a), $D_{v} f \in \mathcal{A}_{\tau_{v}}$. Thus the main point of our lemma is the $\overline{\mathbf{Q}}$ rationality. In Cases SP and UT we can put $f=g_{1} / g_{0}$ with $g_{0}, g_{1} \in \mathcal{M}_{\mu \mathbf{a}}(\overline{\mathbf{Q}}), 0<$ $\mu \in \mathbf{Z}$, since $\mathcal{A}_{0}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ is generated by the quotients of (7.8). Let $z_{j k}^{v}$ be the $(j, k)$ entry of $z_{v}$. Then $\pi^{-1} g_{0}^{2} \partial f / \partial z_{j k}^{v}$ has a Fourier expansion with coefficients in $\overline{\mathbf{Q}}$, as we already observed in (9.9) in a similar setting. This means that $\pi^{-1} g_{0}^{2} D_{v} f \in$ $\mathcal{M}_{\xi}(\overline{\mathbf{Q}})$ with $\xi(a, b)=\operatorname{det}(b)^{2 \mu \mathbf{a}} \tau_{v}(a, b)$. Thus $\pi^{-1} D_{v} f \in \mathcal{A}_{\tau_{v}}(\overline{\mathbf{Q}})$.

In Case UB the matter is not so simple. We employ the notation of $\S 11.6$; in particular we consider the embedding $\varepsilon: \mathcal{H} \rightarrow \mathfrak{H}_{d}$ of (11.8). Also, let $\Gamma,\left(A_{z}, \mathcal{C}_{z}\right)$, and $\mathfrak{F}_{z}$ be as in §11.1. Take a point $z_{0}$ of $\mathcal{H}$ generic for the elements of $\mathcal{A}_{0}(\Gamma, \overline{\mathbf{Q}})$ over $\overline{\mathbf{Q}}$ in the sense of Section 7. Then $\mathfrak{F}_{z_{0}}$ is algebraic over the field of moduli of $\left(A_{z_{0}}, \mathcal{C}_{z_{0}}\right)$, which can be generated by $\mathfrak{f}\left(\varepsilon\left(z_{0}\right)\right)$ for all $\mathfrak{f} \in \mathfrak{A}_{0}\left(\Gamma^{1}, \mathbf{Q}\right)$ finite at $\varepsilon\left(z_{0}\right)$, where $\Gamma^{1}=S p(d, \mathbf{Z})$. Therefore we can find elements $\mathfrak{f}_{1}, \ldots, \mathfrak{f}_{p}$ of $\mathfrak{A}_{0}\left(\Gamma^{1}, \mathbf{Q}\right)$ such that $\mathfrak{f}_{1} \circ \varepsilon, \ldots, \mathfrak{f}_{p} \circ \varepsilon$ are algebraically independent, where $p$ is the complex dimension of $\mathcal{H}$. Put $g_{j}=\mathfrak{f}_{j} \circ \varepsilon$. By (11.12), $g_{j} \in \mathcal{A}_{0}(\overline{\mathbf{Q}})$. Then $\partial / \partial g_{1}, \ldots, \partial / \partial g_{p}$ are well-defined derivations of $\mathcal{A}_{0}(\overline{\mathbf{Q}})$, and for $f \in \mathcal{A}_{0}(\overline{\mathbf{Q}})$ we have $D_{v} f=\sum_{j=1}^{p}\left(\partial f / \partial g_{j}\right) D_{v} g_{j}$. Since $\partial f / \partial g_{j} \in \mathcal{A}_{0}(\overline{\mathbf{Q}})$, our task is to show that $D_{v} g_{j} \in \mathcal{A}_{\tau_{v}}(\overline{\mathbf{Q}})$, or rather $D_{v}(\mathfrak{f} \circ \varepsilon) \in \mathcal{A}_{\tau_{v}}(\overline{\mathbf{Q}})$ for every $\mathfrak{f} \in \mathfrak{A}_{0}(\overline{\mathbf{Q}})$ for which $\mathfrak{f} \circ \varepsilon$ is meaningful. Fixing $v$, put $m=m_{v}$ and $n=n_{v}$. We may assume that $\mathcal{T}=\operatorname{diag}\left[\zeta_{1}, \ldots, \zeta_{r}\right]$ with $\zeta_{\mu} \in K^{\times}$such that $\zeta_{\mu}^{\rho}=-\zeta_{\mu}$. Our $v$ defines an embedding of $F$ into $\mathbf{R}$, and in $\S 3.5$ we fixed an embedding of $K$ into $\mathbf{C}$ that extends $v$. For $x \in K$ denote by $x_{v}$ the image of $x$ by that embedding. By (4.17) we may assume that $i \zeta_{\mu v}>0$ for $\mu \leq m$ and $i \zeta_{\mu v}<0$ for $\mu>m$. Take real numbers $s_{1}, \ldots, s_{r}$ so that $s_{\mu}^{2}=i \zeta_{\mu v}$ for $\mu \leq m$ and $s_{\mu}^{2}=-i \zeta_{\mu v}$ for $\mu>m$. Then we can take $\operatorname{diag}\left[s_{1}, \ldots, s_{r}\right]$ as $Q_{v}$ of (3.34). Put $[F: \mathbf{Q}]=t$. For each $\mu$ take $\mathbf{Q}$-bases $\left\{a_{\mu j}\right\}_{j=1}^{t}$ and $\left\{a_{\mu j}^{\prime}\right\}_{j=1}^{t}$ of $F$ so that

$$
\begin{equation*}
\operatorname{Tr}_{F / \mathbf{Q}}\left(\zeta_{\mu}^{2} a_{\mu j} a_{\mu k}^{\prime}\right)=\delta_{j k} / 2 \tag{14.5}
\end{equation*}
$$

Let $\left\{\tilde{e}_{\mu}\right\}_{\mu=1}^{r}$ be the standard basis of $K_{r}^{1}$; put $h_{\mu j}=a_{\mu j} \tilde{e}_{\mu}$ and $h_{\mu j}^{\prime}=\zeta_{\mu} a_{\mu j}^{\prime} \tilde{e}_{\mu}$. Then the elements

$$
h_{11}, \ldots, h_{1 t}, \cdots \cdots, h_{r 1}, \ldots, h_{r t}, h_{11}^{\prime}, \ldots, h_{1 t}^{\prime}, \cdots \cdots, h_{r 1}^{\prime}, \ldots, h_{r t}^{\prime}
$$

form a $\mathbf{Q}$-basis of $K_{r}^{1}$. Define $g: K_{r}^{1} \rightarrow \mathbf{Q}_{2 d}^{1}$ so that the image of these by $g$ is the standard basis $\left\{e_{k}\right\}_{k=1}^{2 d}$ of $\mathbf{Q}_{2 d}^{1}$. Then (11.5) is satisfied. Let $p_{z}$ and $\kappa(z)$ be as in (11.9) and (11.7). From (4.19) and (4.22) we see that

$$
\left[\begin{array}{lll}
p_{z}\left(\widetilde{e}_{1}\right)_{v} & \cdots & p_{z}\left(\widetilde{e}_{r}\right)_{v}
\end{array}\right]=\left[\begin{array}{cc}
b_{v} & z_{v} c_{v} \\
t_{v} b_{v} & c_{v}
\end{array}\right],
$$

where $b=\operatorname{diag}\left[s_{1}, \ldots, s_{m}\right]$ and $c=\operatorname{diag}\left[s_{m+1}, \ldots, s_{r}\right]$. Now we have $\mathbf{C}^{d}=\left(\mathbf{C}^{r}\right)^{\mathbf{a}}$, and we identify a with $\{1, \ldots, t\}$ so that our fixed $v$ corresponds to the index 1 . Focusing our attention on the first $r$ components of the vectors $p_{z}\left(g^{-1}\left(e_{k}\right)\right)$ and employing (4.10), (4.12), (4.18), (11.6), and (11.7), we see that

$$
\kappa(z) \varepsilon(z)=\left[\begin{array}{cc}
\beta & z_{v} \gamma \\
{ }^{t} z_{v} \beta & \gamma \\
* & *
\end{array}\right], \quad \kappa(z)=\left[\begin{array}{cc}
\beta^{\prime} & z_{v} \gamma^{\prime} \\
-{ }^{t} z_{v} \beta^{\prime} & -\gamma^{\prime} \\
* & *
\end{array}\right]
$$

with ( $m \times m t$ )-matrices $\beta, \beta^{\prime}$ and $(n \times n t)$-matrices $\gamma, \gamma^{\prime}$ whose entries can be given explicitly in terms of the conjugates of $s_{\mu}, a_{\mu j}, a_{\mu j}^{\prime}$, and $\zeta_{\mu}$. Now for a function $\psi$ on $\mathcal{H}$ denote by $d_{v} \psi$ the holomorphic $v$-part of the 1 -form $d \psi$; in other words, $d_{v} \psi=\left(D_{v} \psi\right)\left(d z_{v}\right)$ (see (12.13)). Then

$$
\begin{aligned}
\kappa \cdot d_{v} \varepsilon \cdot{ }^{t} \kappa & =d_{v}(\kappa \varepsilon) \cdot{ }^{t} \kappa-d_{v} \kappa \cdot{ }^{t}(\kappa \varepsilon) \\
& =\left[\begin{array}{ccc}
d z_{v}\left(\gamma \cdot{ }^{t} \gamma^{\prime}-\gamma^{\prime} \cdot{ }^{t} \gamma\right) \cdot{ }^{t} z_{v} & -d z_{v}\left(\gamma \cdot{ }^{t} \gamma^{\prime}+\gamma^{\prime} \cdot{ }^{t} \gamma\right) & * \\
{ }^{t} d z_{v}\left(\beta \cdot{ }^{t} \beta^{\prime}+\beta^{\prime} \cdot{ }^{t} \beta\right) & { }^{t} d z_{v}\left(\beta^{\prime} \cdot{ }^{t} \beta-\beta \cdot{ }^{t} \beta^{\prime}\right) z_{v} & * \\
0 & 0 & 0
\end{array}\right] .
\end{aligned}
$$

Since this matrix is symmetric, using the explicit forms of $\beta, \beta^{\prime}, \gamma, \gamma^{\prime}$ and (14.5), we find that

$$
\kappa \cdot d_{v} \varepsilon \cdot{ }^{t} \kappa=i\left[\begin{array}{ccc}
0 & d z_{v} & 0 \\
{ }^{t} d z_{v} & 0 & 0 \\
0 & 0 & 0
\end{array}\right] .
$$

Defining $D$ on $\mathfrak{H}_{d}$ by $d \mathfrak{f}=(D \mathfrak{f})(d Z)$ for the variable $Z$ on $\mathfrak{H}_{d}$, we have $d_{v}(\mathfrak{f} \circ$ $\varepsilon)=((D \mathfrak{f}) \circ \varepsilon)\left(d_{v} \varepsilon\right)$. Put $\mathfrak{T}=\left\{\left.U \in \mathbf{C}_{d}^{d}\right|^{t} U=U\right\}$ and view it as the (holomorphic) tangent space of $\mathfrak{H}_{d}$. Define $S_{1}(\mathfrak{T})$-valued function $Y$ on $\mathfrak{H}_{d}$ by $Y(U)=$ $\pi^{-1}(D \mathfrak{f})\left({ }^{t} P^{-1} U P^{-1}\right)$ for $U \in \mathfrak{T}$ with $P$ as in $\S 11.8$. Then $Y(U)$ for an algebraic $U$ is an element of $\mathfrak{A}_{0}(\overline{\mathbf{Q}})$. Let $S_{v}, R_{v}$, and $W$ be as in (11.13); put $A(z)(U)=$ $Y(\varepsilon(z))\left({ }^{t} W(z) U W(z)\right)$ and $\mathfrak{R}=\operatorname{diag}\left[S_{v}, R_{v}\right]_{v \in \mathbf{a}}$. Then

$$
\begin{aligned}
& \pi^{-1} D_{v}(f \circ \varepsilon)\left(d z_{v}\right)=\pi^{-1} d_{v}(f \circ \varepsilon)=A\left({ }^{t} W^{-1} \cdot{ }^{t}(P \circ \varepsilon) d_{v} \varepsilon(P \circ \varepsilon) W^{-1}\right) \\
&=A\left({ }^{t} \mathfrak{R} \cdot \kappa \cdot d_{v} \varepsilon \cdot{ }^{t} \kappa \cdot \mathfrak{R}\right)=i A\left(\left[\begin{array}{ccc}
0 & { }^{t} S_{v} d z_{v} R_{v} & 0 \\
{ }^{t} R_{v} \cdot{ }^{t} d z_{v} S_{v} & 0 & 0 \\
0 & 0 & 0
\end{array}\right]\right) .
\end{aligned}
$$

Since $W$ has entries in $\mathcal{A}_{0}(\overline{\mathbf{Q}}), A(w)(U)$ for $w \in \mathcal{H}_{\mathrm{CM}}$ and algebraic $U$ is algebraic. Put $B(z)(u)=A(z)\left(\left[\begin{array}{ccc}0 & u & 0 \\ t_{u} & 0 & 0 \\ 0 & 0 & 0\end{array}\right]\right)$ for $u \in T_{v}$. Then $\pi^{-1} D_{v}(f \circ \varepsilon)(u)=$ $i B\left({ }^{t} S_{v} u R_{v}\right)$. Thus for $w \in \mathcal{H}_{\mathrm{CM}}$ we have
$\left(^{*}\right) \quad \pi^{-1}\left(\mathfrak{P}_{\tau_{v}}(w)^{-1} D_{v}(\mathfrak{f} \circ \varepsilon)(w)\right)(u)=i B(w)\left({ }^{t} S_{v}(w) \cdot{ }^{t} \mathfrak{p}_{v \rho}(w)^{-1} u \mathfrak{p}_{v}(w)^{-1} R_{v}(w)\right)$.
We saw in (11.16) that $\mathfrak{p}_{v}(w)^{-1} R_{v}(w)$ and $\mathfrak{p}_{v \rho}(w)^{-1} S_{v}(w)$ are algebraic for every $w \in \mathcal{H}_{\mathrm{CM}}$ where $P \circ \varepsilon$ and $W$ are finite and invertible. This means that the quantity of $\left(^{*}\right)$ for algebraic $u$ is algebraic. This proves that $\pi^{-1} D_{v}(\mathfrak{f} \circ \varepsilon) \in \mathcal{A}_{\tau_{v}}(\overline{\mathbf{Q}})$ as expected.
14.6. Lemma. Let $\{\omega, X\}$ and $\{\zeta, Y\}$ be $\overline{\mathbf{Q}}$-rational representations of $\mathfrak{K}$, and let $\varphi$ be a $\mathbf{C}$-linear map of $Y$ into $X$ such that $\varphi \zeta(\alpha)=\omega(\alpha) \varphi$ for every $\alpha \in \mathfrak{K}$. Define a map $\varphi_{e}: M l_{e}(T, Y) \rightarrow M l_{e}(T, X)$ by $\varphi_{e}(h)=\varphi \circ h$ for $h \in M l_{e}(T, Y)$. Then the following assertions hold:
(i) $\varphi_{e} \circ\left(\zeta \otimes \tau^{e}\right)(\alpha)=\left(\omega \otimes \tau^{e}\right)(\alpha) \circ \varphi_{e}$ and $\varphi_{e} \circ\left(\zeta \otimes \sigma^{e}\right)(\alpha)=\left(\omega \otimes \sigma^{e}\right)(\alpha) \circ \varphi_{e}$ for every $\alpha \in \mathfrak{K}$.
(ii) $\varphi_{e} C^{e} g=C^{e}(\varphi g), \varphi_{e} E^{e} g=E^{e}(\varphi g)$, and $\varphi_{e} D_{\zeta}^{e} g=D_{\omega}^{e}(\varphi g)$ for every $g \in$ $C^{\infty}(\mathcal{H}, Y)$.
(iii) $\varphi C_{\zeta} \subset C_{\omega}, \varphi \mathcal{M}_{\zeta} \subset \mathcal{M}_{\omega}$, and $\varphi \mathcal{A}_{\zeta} \subset \mathcal{A}_{\omega}$.
(iv) If $\varphi$ is $\overline{\mathbf{Q}}$-rational, $f \in C_{\zeta}$, and $f$ is $\zeta$-arithmetic at $w$, then $\varphi f$ is $\omega$-arithmetic at $w$.

Proof. Assertion (i) follows from (13.19a, b), and the first two equalities of (ii) from (13.10). Employing these and (12.17) (or (13.21)), we obtain the last equality of (ii). Assertions (iii) and (iv) follow immediately from our definition.
14.7. Theorem. Let $Z$ be a $\mathfrak{K}$-stable subspace of $S_{e}(T)$. If $f \in \mathcal{A}_{\omega}(\overline{\mathbf{Q}})$, then $\pi^{-|e|} D_{\omega}^{Z} f$ is arithmetic at every CM-point where $f$ is finite.

Proof. Taking $\varphi_{Z}$ as $\varphi$ of Lemma 14.6, we see that it is sufficient to prove the case $Z=S_{e}(T)$, that is, the case $D_{\omega}^{Z} f=D_{\omega}^{e} f$. We take $X=\mathbf{C}^{m}$. Let $g$ be an $m \times m$-matrix whose columns are all equal to $f$. Define a representation $\{\psi, \operatorname{End}(X)\}$ of $\mathfrak{K}$ by $\psi(x) y=\omega(x) y$ for $y \in \operatorname{End}(X)$. Then our question can be reduced to the nature of $D_{\psi}^{Z} g$. Now given $w \in \mathcal{H}_{\mathrm{CM}}$ where $f$ is finite, take $R=\left(R_{v}\right)_{v \in \mathbf{a}}$ as in Proposition 11.14 (resp. Proposition 9.11) in Case UB (resp, Cases SP and UT), so that the $R_{v}$ are finite and invertible at $w$; here we take the columns of $R_{v}$ in $\mathcal{A}_{\sigma_{v}}(\overline{\mathbf{Q}})$ with $\sigma_{v}(x)=x_{v}$. Put $A_{1}=\omega(R)$ and $A_{2}=c A_{1}+g$ with $c \in \mathbf{Q}$. Then $A_{1}$ is finite and invertible at $w$; we can also take $c$ so that $A_{2}$ is invertible at $w$. Since $g=A_{2}-c A_{1}$, it is sufficient to prove the arithmeticity of $D_{\psi}^{Z} A$ at $w$ for every $A \in \mathcal{A}_{\psi}(\overline{\mathbf{Q}})$ that is finite and invertible at $w$. Fix such an $A$ and take $(Y, h)$ as in $\S 4.11$ such that $w$ is the fixed point of $h\left(Y^{u}\right)$; take also $\beta \in Y^{u}$ as in Lemma 4.12 and put $\gamma=h(\beta)$. Now our method of proof is an adaptation of that of $\S 14.0$. Put $B=A^{-1}\left(A \|_{\psi} \gamma\right)$. Then $A \|_{\psi} \gamma=A B$ and $B$ has entries in $\mathcal{A}_{0}(\overline{\mathbf{Q}})$, and so $B(w)$ is algebraic. Since $\gamma w=w$, we have

$$
\begin{equation*}
B(w)=A(w)^{-1} \omega\left(M_{\gamma}(w)^{-1}\right) A(w) \tag{14.6}
\end{equation*}
$$

By (12.21) we have $\left(D_{\omega}^{e} A\right) \|_{\psi \otimes \tau^{e}} \gamma=D_{\omega}^{e}\left(A \|_{\psi} \gamma\right)=D_{\omega}^{e}(A B)$, and hence, by Lemma 13.17, for $u \in T$ we have

$$
\begin{equation*}
\omega\left(M_{\gamma}^{-1}\right)\left(D_{\omega}^{e} A\right)\left({ }^{t} \lambda_{\gamma}^{-1} u \mu_{\gamma}^{-1}\right)=\sum_{a+b=e}\binom{e}{a}\left(D_{\omega}^{a} A\right)(u)\left(D_{\iota}^{b} B\right)(u) \tag{14.7}
\end{equation*}
$$

where $\iota$ denotes the trivial representation of $\mathfrak{K}$. Put

$$
R_{a}(u)=\pi^{-|a|} A(w)^{-1}\left(D_{\omega}^{a} A\right)(w)\left(\mathfrak{p}_{v \rho}^{-1} u \mathfrak{p}_{v}^{-1}\right), \quad S_{b}(u)=\pi^{-|b|}\left(D_{\iota}^{b} B\right)(w)\left({ }^{t} \mathfrak{p}_{v \rho}^{-1} u \mathfrak{p}_{v}^{-1}\right)
$$

where ${ }^{t} \mathfrak{p}_{v \rho}^{-1} u \mathfrak{p}_{v}^{-1}=\left({ }^{t} \mathfrak{p}_{v \rho}(w)^{-1} u_{v} \mathfrak{p}_{v}(w)^{-1}\right)_{v \in \mathbf{a}}$ for $u=\left(u_{v}\right)_{v \in \mathbf{a}} \in T$. Then $B(w)^{-1}$ times (14.7) gives

$$
\begin{equation*}
R_{e}\left({ }^{t} \zeta^{\prime} u \zeta\right)-B(w)^{-1} R_{e}(u) B(w)=\sum_{a+b=e, a \neq e}\binom{e}{a} B(w)^{-1} R_{a}(u) S_{b}(u) \tag{14.8}
\end{equation*}
$$

Here $\zeta=\left(\mathfrak{p}_{v}(w)^{-1} \mu_{v}(\gamma, w)^{-1} \mathfrak{p}_{v}(w)\right)_{v \in \mathbf{a}}$ and $\zeta^{\prime}=\left(\mathfrak{p}_{v \rho}(w)^{-1} \lambda_{v}(\gamma, w)^{-1} \mathfrak{p}_{v \rho}(w)\right)_{v \in \mathbf{a}}$. Since $A \in \mathcal{A}_{\psi}(\overline{\mathbf{Q}})$, we see that $\mathfrak{P}_{\omega}(w)^{-1} A(w)$ is algebraic. Therefore our task is to prove the algebraicity of $R_{e}(u)$ for algebraic $u$. We are going to prove this by induction on $|e|$. Since $\gamma w=w$, Proposition 11.5 (2) shows that $\zeta$ and $\zeta^{\prime}$ are algebraic. Also, the linear endomorphism $R(u) \mapsto R\left({ }^{t} \zeta^{\prime} u \zeta\right)$ of $S_{e}(T, \operatorname{End}(X))$ commutes with another endomorphism $R \mapsto B(w)^{-1} R B(w)$. Now we can choose $\gamma$ so that $\left(\zeta^{\prime} \otimes \zeta\right)^{\otimes|e|}$ and $B(w)^{-1} \otimes B(W)$ have no common eigenvalue, as will be shown at the end of the proof. Take any algebraic $u$. By the induction alssumption, $R_{a}(u)$ on the right-hand side of (14.8) is algebraic. As for $S_{b}(u)$, we have $b \neq 0$, so that
we have $D_{\iota}^{b} B=D_{\tau_{v}}^{c}\left(D_{v} B\right)$ for some $c$ and some $v \in \mathbf{a}$. By Proposition 14.5, the entries of $\pi^{-1} D_{v} B$ belong to $\mathcal{A}_{\tau_{v}}(\overline{\mathbf{Q}})$. Since $|c|<|e|$, the induction assumption implies the algebraicity of $S_{b}(u)$. Thus the sum on the right-hand side of (14.8) is algebraic. Viewing (14.8) as a system of linear equations with $R_{e}(u)$ as the unknown, we obtain the desired algebraicity of $R_{e}(u)$. Notice that if $|e|=1$, the right-hand side of (14.8) consists of a single term $\pi^{-1} B(w)^{-1}\left(D_{v} B\right)(w)\left({ }^{t} \mathfrak{p}_{v \rho}(w)^{-1} u \mathfrak{p}_{v}(w)^{-1}\right)$ with $u \in T_{v}$ for some $v$. Therefore Proposition 14.5 gives its algebraicity.

To prove the existence of $\gamma=h(\beta)$ with the desired property on eigenvalues, consider $\Phi$ and $\left(K_{i}, \Phi_{i}\right)$ defined for the present $Y$ as in $\S 4.11 ; \operatorname{put}[Y: \mathbf{Q}]=2 d$. Then there exist $d$ maps $y \mapsto y_{i}$ for $1 \leq i \leq d$ of $Y$ into $\mathbf{C}$ such that $y_{1}, \ldots, y_{d}$ are the eigenvalues of $\Phi(y)$; moreover these maps together with their complex conjugates form exactly the set $J_{Y}$ of $\S 11.3$. By (4.37) and (4.40) we see that $\alpha_{1}, \ldots, \alpha_{d}$ are exactly the eigenvalues of $M(h(\alpha), w)$ for every $\alpha \in Y^{u}$. Clearly we may assume that $\omega$ is irreducible. Then we can find a homogeneous polynomial representation $\sigma$ of $\mathfrak{K}$ such that $\sigma(a, b)=\operatorname{det}(a)^{\mu} \operatorname{det}(b)^{\nu} \omega(a, b)$ with $\mu, \nu \in \mathbf{Z}^{\mathrm{a}}$. From (14.6) we see that $B(w)^{-1} \otimes B(w)$ has the same eigenvalues as $\sigma\left(M_{\gamma}(w)\right)^{-1} \otimes \sigma\left(M_{\gamma}(w)\right)$. Therefore each eigenvalue of $B(w)^{-1} \otimes B(w)$ is of the form $\prod_{i=1}^{d} \beta_{i}^{\lambda_{i}}$ with integers $\lambda_{i}$ such that $\sum_{i=1}^{d} \lambda_{i}=0$. On the other hand each eigenvalue of $\left(\zeta^{\prime} \otimes \zeta\right)^{\otimes|e|}$ is of the form $\prod_{i=1}^{d} \beta_{i}^{-\mu_{i}}$ with nonnegative integers $\mu_{i}$ such that $\sum_{i=1}^{d} \mu_{i}=2|e|$. (In Case SP we have $\lambda_{v}(\gamma, w)=\mu_{v}(\gamma, w)$ and $\zeta^{\prime}=\zeta$, but still our statements are valid.) Suppose that $B(w)^{-1} \otimes B(w)$ and $\left(\zeta^{\prime} \otimes \zeta\right)^{\otimes|e|}$ have a common eigenvalue. Then there exist $d$ integers $\kappa_{i}$ such that $\prod_{i=1}^{d} \beta_{i}^{\kappa_{i}}=1, \sum_{i=1}^{d} \kappa_{i}=2|e|$, and $\left|\kappa_{i}\right| \leq N$ with a positive integer $N$ depending on $\omega$ and $e$. Now the map $\alpha \mapsto\left(\alpha_{1}, \ldots, \alpha_{d}\right)$ sends $Y^{u}$ into a dense subset of $\mathbf{T}^{d}$, since $Y$ is dense in $\mathbf{C}^{d}$ and $x^{\rho} / x \in Y^{u}$ for every $x \in Y^{\times}$. Therefore we can find an element $\beta$ of $Y^{u}$ such that $\prod_{i=1}^{d} \beta_{i}^{\kappa_{i}} \neq 1$ for every $\left\{\kappa_{i}\right\}$ as above. This proves the desired fact, and our proof is now complete.
14.8. Proposition. Let $(V, \varphi)$ be a model of $\Gamma / \mathcal{H}$ with the properties given in Theorem 9.1 in Case UB; suppose that $\Gamma \subset S U(\mathcal{T})$. Then the following assertions hold:
(1) If $0 \neq h \in \mathcal{A}_{\kappa \mathbf{b}}(\Gamma, \overline{\mathbf{Q}})$ with $\kappa \in \mathbf{Z}$, then $\operatorname{div}(h)$ considered on $V$ is $\overline{\mathbf{Q}}$-rational.
(2) There exist a positive integer $m$ and a nonzero element $g \in \mathcal{A}_{m \mathbf{b}}(\Gamma, \overline{\mathbf{Q}})$ such that $\operatorname{div}(g)$ considered on $V$ is $\overline{\mathbf{Q}}$-rational.

Proof. Let $p$ be the complex dimension of $\mathcal{H}$. Take $p$ algebraically independent functions $h_{1}, \ldots, h_{p}$ in $\mathcal{A}_{0}(\Gamma, \overline{\mathbf{Q}})$. Let $z_{a b}^{v}$ be the ( $a, b$ )-entry of the matrix $z_{v}$ which is the $v$-th component of the variable $z=\left(z_{v}\right)_{v \in \mathbf{a}} \in \mathcal{H}$. Put

$$
q=\pi^{-p} \partial\left(h_{1}, \ldots, h_{p}\right) / \partial\left(z_{1}, \ldots, z_{p}\right)
$$

where $z_{1}, \ldots, z_{p}$ are an arbitrarily fixed arrangement of the variables $z_{a b}^{v}$ for all $v \in \mathbf{a}^{\prime}$ and all $(a, b)$. From Lemma 3.4 (2) we see that $q^{2} \in \mathcal{A}_{\zeta}(\Gamma)$, where $\zeta(x)=$ $\operatorname{det}(x)^{r \mathbf{b}}$. Now we consider $D_{v} h_{i}$ for $v \in \mathbf{a}^{\prime}$. By Proposition 14.5, $\pi^{-1}\left(D_{v} h_{i}\right)(w)$ $\left({ }^{t_{p}}{ }_{v \rho}(w)^{-1} u \mathfrak{p}_{v}(w)^{-1}\right)$ is $\overline{\mathbf{Q}}$-rational for every $\overline{\mathbf{Q}}$-rational $u \in T_{v}$ and for every $w \in \mathcal{H}_{\mathrm{CM}}$ where $h_{i}$ is finite. From this we see that $q^{2} \in \mathcal{A}_{\xi}(\Gamma, \overline{\mathbf{Q}})$ with $\xi(x)=$ $\prod_{v \in \mathbf{a}^{\prime}} \operatorname{det}\left(x_{v \rho}\right)^{2 n_{v}} \operatorname{det}\left(x_{v}\right)^{2 m_{v}}$. By Proposition $11.17(2), \mathcal{A}_{\xi}(\overline{\mathbf{Q}})=\mathfrak{q} \cdot \mathcal{A}_{\zeta}(\overline{\mathbf{Q}})$ with a certain constant $\mathfrak{q}$. Thus $\mathfrak{q}^{-1} q^{2} \in \mathcal{A}_{r \mathbf{b}}(\Gamma, \overline{\mathbf{Q}})$. Now $\operatorname{div}(q)$ considered on $V$ is the divisor of $d h_{1} \wedge \cdots \wedge d h_{p}$, which is $\overline{\mathbf{Q}}$-rational. Taking $g=\mathfrak{q}^{-1} q^{2}$, we obtain (2). Then assertion (1) can be proved in exactly the same fashion as in the proof of Proposition 9.8 (1).

Proof of (11.19). Let $h=\prod_{v \in \mathbf{b}} \operatorname{det}\left(R_{v}\right)^{\kappa}$ with $R_{v}$ as in Proposition 11.14. Then $0 \neq h \in \mathcal{A}_{\kappa \mathbf{b}}(\Gamma, \overline{\mathbf{Q}})$ with a suitable $\Gamma$ as in Proposition 14.8. Then (11.19) can be proved by the same technique as in the proof of Theorem 9.9 (1), (2).
14.9. Theorem. (1) Let $\mathcal{W}$ be a dense subset of $\mathcal{H}$ contained in $\mathcal{H}_{\mathrm{CM}}$ (see $\S 11.1)$, and $\{\omega, X\}$ a $\overline{\mathbb{Q}}$-rational representation of $\mathfrak{K}$. If an element $f$ of $\mathcal{N}_{\omega}^{p}$ is $\omega$-arithmetic at every point of $\mathcal{W}$, then $f \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$.
(2) Let $Z$ be an irreducible subspace of $S_{e}(T)$. Then $\pi^{-|e|} D_{\omega}^{Z} f \in \mathcal{N}_{\omega \otimes \tau_{Z}}^{p+e}(\overline{\mathbf{Q}})$ and $\pi^{|e|} E^{Z} f \in \mathcal{N}_{\omega \otimes \sigma_{Z}}^{p^{\prime}}(\overline{\mathbf{Q}})$ for every $f \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$, where $|e|=\sum_{v \in \mathbf{a}} e_{v}$ and $p_{v}^{\prime}=$ $\operatorname{Max}\left(p_{v}-e_{v}, 0\right)$.

Proof. Fix $p$; let $\mathcal{N}_{\omega}^{\prime}$ be the set of all $f \in \mathcal{N}_{\omega}^{p}$ that are $\omega$-arithmetic at every point of $\mathcal{W}$. Put $\zeta(\alpha)=\operatorname{det}(\alpha)^{\kappa \mathbf{b}} \omega(\alpha)$ for $\alpha \in \mathfrak{K}$ with $0<\kappa \in \mathbf{Z}$. Let $\mathcal{N}_{\zeta}^{*}$ (resp. $\left.\mathcal{N}_{\zeta}^{*}(\overline{\mathbf{Q}})\right)$ be the set of all $h \in \mathcal{N}_{\zeta}^{p}$ of the form

$$
\begin{equation*}
h=\sum_{s \leq p} \pi^{-|s|} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{s} \tag{14.9}
\end{equation*}
$$

with $g_{s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}\left(\right.$ resp. $g_{s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\overline{\mathbf{Q}})$ ). By Proposition 14.2, $\mathcal{N}_{\zeta}^{p}=\mathcal{N}_{\zeta}^{*}$ if $\kappa$ is sufficiently large for every $v \in \mathbf{a}^{\prime}$. Fix such a $\kappa$. By Proposition 11.15 we see that $\mathcal{N}_{\zeta}^{p}$ is spanned by $\mathcal{N}_{\zeta}^{*}(\overline{\mathbf{Q}})$ over $\mathbf{C}$. By Theorem 14.7 and Lemma 14.6 (iv), $\pi^{-|s|} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{s}$ is $\zeta$-arithmetic at every point of $\mathcal{H}_{\mathrm{CM}}$ for every $g_{s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\overline{\mathbf{Q}})$, and hence $\mathcal{N}_{\zeta}^{*}(\overline{\mathbf{Q}}) \subset \mathcal{N}_{\zeta}^{p}(\overline{\mathbf{Q}}) \subset \mathcal{N}_{\zeta}^{\prime}$. Let us now prove that $\mathcal{N}_{\zeta}^{\prime} \subset \mathcal{N}_{\zeta}^{*}(\overline{\mathbf{Q}})$. Take a basis $B$ of $\mathbf{C}$ over $\overline{\mathbf{Q}}$ including 1 ; let $h \in \mathcal{N}_{\zeta}^{\prime}$. Then $h=\sum_{c \in C} c k_{c}$ with a finite subset $C$ of $B$ and $k_{c} \in \mathcal{N}_{\zeta}^{*}(\overline{\mathbf{Q}})$. For every $w \in \mathcal{W}$ we have $\mathfrak{P}_{\zeta}(w)^{-1} h(w)=\sum_{c} c \mathfrak{P}_{\zeta}(w)^{-1} k_{c}(w)$. Since $\mathfrak{P}_{\zeta}(w)^{-1} h(w)$ and $\mathfrak{P}_{\zeta}(w)^{-1} k_{c}(w)$ are algebraic, we have $k_{c}(w)=0$ for $c \neq 1$, and hence $h(w)=k_{1}(w)$ for every $w \in \mathcal{W}$. Since $\mathcal{W}$ is dense in $\mathcal{H}$, we have $h=k_{1}$. This proves that $\mathcal{N}_{\zeta}^{\prime}=\mathcal{N}_{\zeta}^{*}(\overline{\mathbf{Q}})=\mathcal{N}_{\zeta}^{p}(\overline{\mathbf{Q}})$. To show that $\mathcal{N}_{\omega}^{\prime} \subset \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$, take any $w_{0} \in \mathcal{H}_{\mathrm{CM}}$ and take $q \in \mathcal{M}_{\kappa_{0} \mathbf{b}}(\overline{\mathbf{Q}})$ with some $\kappa_{0} \in \mathbf{Z}$ so that $q\left(w_{0}\right) \neq 0$. Such a $q$ in Cases SP and UT is obtained in Lemma 6.17; in Case UB take $q=\prod_{v \in \mathrm{~b}} \operatorname{det}\left(Q_{v}\right)$ with $Q_{v}$ of Proposition 11.14. Changing $q$ for its suitable power, we may take $\kappa=\kappa_{0}$. Let $f \in \mathcal{N}_{\omega}^{\prime}$. Clearly $q f \in \mathcal{N}_{\zeta}^{\prime}=\mathcal{N}_{\zeta}^{p}(\overline{\mathbf{Q}})$. We can take $q\left(w_{0}\right) \mathfrak{P}_{\omega}\left(w_{0}\right)$ as $\mathfrak{P}_{\zeta}\left(w_{0}\right)$. Then $\mathfrak{P}_{\omega}\left(w_{0}\right)^{-1} f\left(w_{0}\right)=\mathfrak{P}_{\zeta}\left(w_{0}\right)^{-1}(q f)\left(w_{0}\right)$, which is algebraic. Since this is so for every $w_{0} \in \mathcal{H}_{\mathrm{CM}}$, we see that $f \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$. This proves (1).

Next, by Lemma 14.6, in order to prove that $\pi^{-|e|} D_{\omega}^{Z} f \in \mathcal{N}_{\omega \otimes \tau_{Z}}^{p+e}(\overline{\mathbf{Q}})$, it is sufficient to show that $\pi^{-|e|} D_{\omega}^{e} f \in \mathcal{N}_{\omega \otimes \tau^{e}}^{p+e}(\overline{\mathbf{Q}})$. Take $w_{0}$ and $q$ again; then we can put $q f=\sum_{s \leq p} \ell_{s}, \ell_{s}=\pi^{-|s|} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{s}$ with $g_{s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\overline{\mathbf{Q}})$. Thus our question is the arithmeticity of $\pi^{-|e|} D_{\omega}^{e}\left(q^{-1} \ell_{s}\right)$. First, Lemma 13.17 shows that

$$
D_{\omega}^{e}\left(q^{-1} \ell_{s}\right)(u)=\sum_{a+b=e}\binom{e}{a}\left(D_{-\kappa}^{a}\left(q^{-1}\right)\right)(u)\left(D_{\zeta}^{b} \ell_{s}\right)(u)
$$

for $u \in T$, where $-\kappa$ stands for $-\kappa \mathbf{b}$. Put $\mathfrak{p}_{1}=\left(\mathfrak{p}_{v \rho}\left(w_{0}\right)\right)_{v \in \mathbf{a}}$ and $\mathfrak{p}_{2}=\left(\mathfrak{p}_{v}\left(w_{0}\right)\right)_{v \in \mathbf{a}}$. Then we can write $\mathfrak{p}\left(w_{0}\right)=\left(\mathfrak{p}_{1}, \mathfrak{p}_{2}\right)$. For $\psi \in S_{e}(T, X)$ we have $\left[\mathfrak{P}_{\omega \otimes \tau^{e}}\left(w_{0}\right)^{-1} \psi\right](u)$ $=\mathfrak{P}_{\omega}\left(w_{0}\right)^{-1} \psi\left({ }^{t} \mathfrak{p}_{1}^{-1} u \mathfrak{p}_{2}\right)$. Now $\pi^{-|a|} D_{-\kappa}^{a}\left(q^{-1}\right)$ is arithmetic by Theorem 14.7. On the other hand, by Lemma 14.6 (ii) we have $\pi^{-|b|} D_{\zeta^{b}} \ell_{s}=\pi^{-|s|-|b|} D_{\zeta}^{b} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{s}=$ $\pi^{-|s|-|b|}\left(\theta_{X}^{s}\right)_{b} D_{\zeta \otimes \sigma^{s}}^{s+b} g_{s}$. Hence by Theorem 14.7 and Lemma 14.6 (i, iv) $\pi^{-|b|} D_{\zeta}^{b} \ell_{s}$ is arithmetic. Thus
$\mathfrak{P}_{\omega \otimes \tau^{e}}\left(w_{0}\right)^{-1} \pi^{-|e|} D_{\omega}^{e}\left(q^{-1} \ell_{s}\right)(u)$

$$
=\sum_{a, b}\binom{e}{a} q\left(w_{0}\right)\left[\pi^{-|a|}\left(D_{-\kappa}^{a}\left(q^{-1}\right)\right)\right]\left({ }^{t} \mathfrak{p}_{1}^{-1} u \mathfrak{p}_{2}\right) \mathfrak{P}_{\zeta}\left(w_{0}\right)^{-1}\left[\pi^{-|b|}\left(D_{\zeta}^{b} \ell_{s}\right)\right]\left({ }^{t} \mathfrak{p}_{1}^{-1} u \mathfrak{p}_{2}\right)
$$

which is algebraic for every algebraic $u \in T$. This completes the proof of the part of (2) concerning $D_{\omega}^{Z}$.

Finally, to deal with $E^{Z}$, it is sufficient to show that $\pi^{|e|} E^{e} f \in \mathcal{N}_{\omega \otimes \sigma^{e}}^{p^{\prime}}(\overline{\mathbf{Q}})$ if $f \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$. Clearly it is sufficient to prove it when $|e|=1$, that is, $E^{e}=E_{v}$ with some $v$. We prove this by induction on $|p|$. If $p=0$, then $f$ is holomorphic and $E_{v} f=0$. Thus we take $|p|>0$. First we assume that $f=\pi^{-|p|} D_{\psi}^{p} g$ with some $g \in \mathcal{M}_{\psi}(\overline{\mathbf{Q}})$ and $\overline{\mathbf{Q}}$-trational $\psi$. Then $D_{\psi}^{p} g=D_{\psi^{\prime}, v^{\prime}} D_{\psi}^{p-v^{\prime}} g$ with some $\psi^{\prime}$ and $v^{\prime}$. If $v \neq v^{\prime}$, then $\pi E_{v} f=\pi^{-1} D_{\psi^{\prime \prime}, v^{\prime}} \pi^{2-|p|} E_{v} D_{\psi}^{p-v^{\prime}} g$ with some $\psi^{\prime \prime}$. Now $\pi^{1-|p|} D_{\psi}^{p-v^{\prime}} g \in \mathcal{N}_{\psi^{\prime}}^{p-v^{\prime}}(\overline{\mathbf{Q}})$ by Theorem 14.7, and so $\pi^{2-|p|} E_{v} D_{\psi}^{p-v^{\prime}} g$ is arithmetic by our induction. Then, by what we already proved, $\pi^{-1} D_{\psi^{\prime \prime}, v^{\prime}} \pi^{2-|p|} E_{v} D_{\psi}^{p-v^{\prime}} g$ is arithmetic. Next, assume $v=v^{\prime}$. Then, by Lemma 13.16,

$$
\pi E_{v} f=\pi^{1-|p|} E_{v} D_{\psi^{\prime}, v} D_{\psi}^{b} g=\pi^{1-|p|} D_{\psi^{\prime} \otimes \sigma_{v}, v} E_{v} D_{\psi}^{b} g-A \pi^{-|b|} D_{\psi}^{b} g
$$

with $b=p-v$ and $A \in S_{1}\left(T_{v}, S_{1}\left(T_{v}, \operatorname{End}(Y)\right)\right)$ defined by $A\left(u, u^{\prime}\right)=P_{\psi \otimes \tau^{b}}^{v}\left(u^{\prime}, u\right)$ for $u, u^{\prime} \in T_{v}$, where $Y$ is the representation space of $\psi \otimes \tau^{b}$. For the same reason as in the case $v \neq v^{\prime}$, the first term on the right-hand side and $\pi^{-|b|} D_{\psi}^{b} g$ are arithmetic. To deal with the whole second term, put $h=\pi^{-|b|} D_{\psi}^{b} g$ and $\varphi=$ $\psi \otimes \tau^{b}$; let $w \in \mathcal{H}_{\mathrm{CM}}$. Then $\left[\mathfrak{P}_{\varphi \otimes \tau_{v} \otimes \sigma_{v}}(w)^{-1} A h\right]\left(u, u^{\prime}\right)=A\left(u, u^{\prime}\right) \mathfrak{P}_{\varphi}(w)^{-1} h$ by Proposition 13.15 (4). Since $A$ is $\overline{\mathbf{Q}}$-rational, we can thus establish the arithmeticity of $\pi E_{v} f$ for $f=\pi^{-|p|} D_{\psi}^{p} g$.

Now take an arbitrary element $f \in \mathcal{N}_{\mathcal{W}}^{p}(\overline{\mathbf{Q}})$. We again employ the expression $q f=\sum_{s} \pi^{-|s|} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{s}$ with $g_{s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\overline{\mathbf{Q}})$. Since $q$ is holomorphic, we have $q \cdot \pi E_{v} f=\pi E_{v}(q f)=\sum_{s} \pi^{1-|s|} E_{v} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{s}$. By what we proved about $E_{v} \pi^{-|p|} D_{\psi}^{p} g$, each term of the last sum is arithmetic. Since $q$ is arithmetic, $\pi E_{v} f$ is arithmetic at every CM-point where $q$ is not zero. By (1) this implies that $E_{v} f$ is arithmetic. This completes the proof.

We note here a simple fact:
(14.9a) The symbol $\zeta$ being as above, let $\mathcal{N}_{\zeta}^{*}(\Gamma, \overline{\mathbf{Q}})$ denote the set of elements $h$ of the form (14.9) with $g_{s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\Gamma, \overline{\mathbf{Q}})$. Then $\mathcal{N}_{\zeta}^{*}(\Gamma, \overline{\mathbf{Q}})=\mathcal{N}_{\zeta}^{p}(\Gamma, \overline{\mathbf{Q}})$.
We have seen that $\mathcal{N}_{\zeta}^{*}(\Gamma, \overline{\mathbf{Q}}) \subset \mathcal{N}_{\zeta}^{p}(\Gamma, \overline{\mathbf{Q}})$. Let $h \in \mathcal{N}_{\zeta}^{p}(\Gamma, \overline{\mathbf{Q}})$. By Proposition 14.2 we have (14.9) with $g_{s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\Gamma)$. Applying Proposition 11.15 to the last set, we have $g_{s}=\sum_{b \in B} b g_{s, b}$ with $g_{s, b} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\Gamma, \overline{\mathbf{Q}})$. In this way we can put $h=\sum_{c \in C} c k_{c}$ with a finite subset $C$ of $B$ and $k_{c} \in \mathcal{N}_{\zeta}^{*}(\Gamma, \overline{\mathbf{Q}})$. We have seen in the above proof that $h=k_{1}$. This proves (14.9a).
14.10. Proposition. $\mathcal{N}_{\omega}^{p}=\mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}}) \otimes_{\overline{\mathbf{Q}}} \mathbf{C}$ and $\mathcal{N}_{\omega}^{p}(\Gamma)=\mathcal{N}_{\omega}^{p}(\Gamma, \overline{\mathbf{Q}}) \otimes_{\overline{\mathbf{Q}}} \mathbf{C}$.

Proof. Let $f_{1}, \ldots, f_{\mu}$ be elements of $\mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$ linearly independent over $\overline{\mathbf{Q}}$. Suppose $\sum_{i=1}^{\mu} a_{i} f_{i}=0$ with $a_{i} \in \mathbf{C}$. Let $\left\{b_{i}\right\}_{i=1}^{m}$ be a $\overline{\mathbf{Q}}$-basis of $\sum_{i=1}^{\mu} \overline{\mathbf{Q}} a_{i}$ and let $a_{i}=\sum_{j} c_{i j} b_{j}$ with $c_{i j} \in \overline{\mathbf{Q}}$. Then $\sum_{j} b_{j} \sum_{i} c_{i j} f_{i}=0$, and so $\sum_{j} b_{j} \sum_{i} c_{i j} \mathfrak{P}_{\omega}(w)^{-1}$ $\cdot f_{i}(w)=0$ for every $w \in \mathcal{H}_{\mathrm{CM}}$. Since $\mathfrak{P}_{\omega}(w)^{-1} f_{i}(w)$ is algebraic, we have $\sum_{i} c_{i j}$ $\cdot f_{i}(w)=0$ for every $j$ and every $w \in \mathcal{H}_{\mathrm{CM}}$. Therefore $\sum_{i} c_{i j} f_{i}=0$ for every $j$,
so that $c_{i j}=0$. Thus $a_{i}=0$, which shows that the $f_{i}$ are linearly independent over $\mathbf{C}$. To prove that $\mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$ spans $\mathcal{N}_{\omega}^{p}$ over $\mathbf{C}$, take $\zeta, B$, and $q$ as in the proof of Theorem 14.9. We have seen that $\mathcal{N}_{\zeta}^{p}(\overline{\mathbf{Q}})$ spans $\mathcal{N}_{\zeta}^{p}$ over $\mathbf{C}$. Let $f \in \mathcal{N}_{\omega}^{p}$; then $q f \in \mathcal{N}_{\zeta}^{p}$, and so $q f=\sum_{c \in B} c t_{c}$ with $t_{c} \in \mathcal{N}_{\zeta}^{p}(\overline{\mathbf{Q}})$. Thus $f=\sum_{c \in B} c t_{c} / q$. If we change $q$ for $q^{\prime}$ and $q^{\prime} f=\sum_{c \in B} c t_{c}^{\prime}$ with $t_{c}^{\prime} \in \mathcal{N}_{\zeta}^{p}(\overline{\mathbf{Q}})$, then $q^{\prime} t_{c}=q t_{c}^{\prime}$. Given an arbitrary point $w_{0}$ of $\mathcal{H}$, we can choose a $\overline{\mathbf{Q}}$-rational $q$ so that $q\left(w_{0}\right) \neq 0$. Then $t_{c} / q$ is nearly holomorphic of degree $\leq p$ in a neighborhood of $w_{0}$. Since $t_{c} / q=t_{c}^{\prime} / q^{\prime}$, this shows that $t_{c} / q \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$. This proves the first equality. Next let $f \in \mathcal{N}_{\omega}^{p}(\Gamma)$. Then $f=\sum_{c \in B} c g_{c}$ with $g_{c} \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$, and $\sum_{c} c\left(g_{c}-g_{c} \| \gamma\right)=$ $f-f \| \gamma=0$ for every $\gamma \in \Gamma$. Since $g_{c} \| \gamma \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$, we obtain $g_{c} \| \gamma=g_{c}$, so that $g_{c} \in \mathcal{N}_{\omega}^{p}(\Gamma, \overline{\mathbf{Q}})$. This proves the second equality.
14.11. We now restrict our treatment to Cases SP and UT. Let $\{\omega, X\}$ be a representation of $\mathfrak{K}$. If $f \in \mathcal{N}_{\omega}^{p}(\Gamma)$, then $f(z)=\sum_{e \leq p} g_{e}(z)(r(z))$ with holomorphic maps $g_{e}: \mathcal{H} \rightarrow S_{e}(T, X)$ as noted in (13.23). Recall that $r_{v}(z)=\left({ }^{t} z_{v}-\bar{z}_{v}\right)^{-1}$. Therefore we easily see that $g_{e}(z+b)=g_{e}(z)$ if $\left[\begin{array}{ll}1 & b \\ 0 & 1\end{array}\right] \in \Gamma$, and hence $g_{e}(z)=$ $\sum_{h \in L} c_{e}(h) \mathbf{e}_{\mathbf{a}}^{n}(h z)$ with $c_{e}(h) \in S_{e}(T, X)$ and a lattice $L$ in the vector space $S$ of (5.17). Thus we can put

$$
\begin{equation*}
f(z)=\sum_{h \in L} s_{h}(r(z)) \mathbf{e}_{\mathbf{a}}^{n}(h z) \tag{14.10}
\end{equation*}
$$

with $s_{h} \in \bigcup_{e \leq p} S_{e}(T, X)$. Now, for $\gamma=\operatorname{diag}[a, \widehat{a}] \in \Gamma$ we have $\omega\left({ }^{t} a, a^{*}\right) f\left(a z a^{*}\right)=$ $f(z)$ by (5.19). Since $r_{v}\left(a z a^{*}\right)={ }^{t} a_{v}^{-1} r_{v}(z) \bar{a}_{v}^{-1}$, we easily see that $g_{e}(z)=$ $\left(\omega \otimes \sigma^{e}\right)\left({ }^{t} a, a^{*}\right) g_{e}\left(a z a^{*}\right)$. Thus $g_{e}$ satisfies (5.21) for a suitable $U$ with $\omega \otimes \sigma^{e}$ in place of $\omega$. Suppose now $n>1$ or $F \neq \mathbf{Q}$. Then from Proposition 5.7 we can conclude that $c_{e}(h) \neq 0$ only if $h_{v} \geq 0$ for every $v \in \mathbf{a}$. Consequently $s_{h}$ of (14.10) is not zero only if $h_{v} \geq 0$ for every $v \in \mathbf{a}$. If $n=1$ and $F=\mathbf{Q}$, we need the cusp condition (13.18a).

To speak of the rationality of $f$ over a number field, we assume that $\{\omega, X\}$ has a $\mathbf{Q}$-rational structure, and write the expansion of (14.10) in the form

$$
\begin{equation*}
f(z)=\sum_{h \in L} q_{h}\left(\pi^{-1} i \cdot r(z)\right) \mathbf{e}_{\mathbf{a}}^{n}(h z) \tag{14.11}
\end{equation*}
$$

with polynomials $q_{h}$ belonging to $\bigcup_{e \leq p} S_{e}(T, X)$. Notice that in Case SP, $\pi^{-1} i r(z)$ $=(2 \pi \cdot \operatorname{Im}(z))^{-1}$. Given $\varepsilon \in \operatorname{Aut}(\mathbf{C})$, for $u \in T=\prod_{v \in \mathbf{a}} T_{v}$ we define $u^{[\varepsilon]} \in T$ by

$$
\left(u^{[\varepsilon]}\right)_{v}=\left\{\begin{align*}
u_{v^{\prime}} & \text { if } v \varepsilon=v^{\prime} \text { on } K,  \tag{14.12}\\
{ }^{t} u_{v^{\prime}} & \text { if } v \varepsilon=v^{\prime} \rho \text { on } K,
\end{align*}\right.
$$

where $\rho$ is complex conjugation. In Case UT we are viewing each $v$ as an embedding of $K$ into $\mathbf{C}$ (see §3.5). For $q \in S_{e}(T, X)$ we define $q^{\varepsilon} \in S_{e}(T, X)$ with respect to the natural $\mathbf{Q}$-rational structure of $S_{e}(T, X)$. Then for $f$ as in (14.11) we define $f^{\varepsilon}$ formally by

$$
\begin{equation*}
f^{\varepsilon}(z)=\sum_{h \in L} q_{h}^{\varepsilon}\left(\pi^{-1} i \cdot r(z)^{[\varepsilon]}\right) \mathbf{e}_{\mathbf{a}}^{n}(h z) \tag{14.13}
\end{equation*}
$$

This includes, as a special case, what we defined in $\S 5.10$. Given a subfield $W$ of $\mathbf{C}$, we say that $f$ is $W$-rational if $q_{h}$ is $W$-rational for every $h$. We denote by $\mathcal{N}_{\omega}^{p}(W)$ the set of all $W$-rational elements of $\mathcal{N}_{\omega}^{p}$, and put $\mathcal{N}_{\omega}^{p}(\Gamma, W)=\mathcal{N}_{\omega}^{p}(\Gamma) \cap \mathcal{N}_{\omega}^{p}(W)$. If
$W=\overline{\mathbf{Q}}$, this is consistent with what was defined in $\S 14.4$; see (1) of the following proposition.

Now for $e \in \mathbf{Z}^{\mathbf{a}}$ we define $e \varepsilon \in \mathbf{Z}^{\mathbf{a}}$ by $(e \varepsilon)_{v^{\prime}}=e_{v}$ if $v^{\prime}=v \varepsilon$ on $F$. (In other words, $e \varepsilon$ is defined so that $a^{e \varepsilon}=\left(a^{e}\right)^{\varepsilon}$ for $a \in F^{\times}$.)
14.12. Theorem (Cases SP and UT). (1) The set $\mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$ defined in $\S 14.4$ consists of all the $\overline{\mathbf{Q}}$-rational elements of $\mathcal{N}_{\omega}^{p}$ in the above sense.
(2) For every $\varepsilon \in \operatorname{Aut}(\mathbf{C})$ and $f \in \mathcal{N}_{\omega}^{p}$ the formal series $f^{\varepsilon}$ defined by (14.13) is indeed an element of $\mathcal{N}_{\omega \in}^{p \varepsilon}$.
(3) We have $\left((\pi i)^{-|e|} D_{\omega}^{e} f\right)^{\varepsilon}\left(u^{[\varepsilon]}\right)=\left((\pi i)^{-|e|} D_{\omega^{e}}^{e \varepsilon} f^{\varepsilon}\right)(u)$ and $\left((\pi i)^{|e|} E^{e} f\right)^{\varepsilon}\left(u^{[\varepsilon]}\right)$ $=\left((\pi i)^{|e|} E^{e \varepsilon} f^{\varepsilon}\right)(u)$ for every $f \in \mathcal{N}_{\omega}^{p}$ and $u \in T$.
(4) Let $Z$ be a $\mathfrak{K}$-stable subspace of $S_{e}(T)$, and $W$ a subfield of $\mathbf{C}$ containing the Galois closure of $K$ over $\mathbf{Q}$. Then $(\pi i)^{-|e|} D_{\omega}^{Z}$ (resp. ( $\left.\pi i\right)^{|e|} E^{Z}$ ) sends $\mathcal{N}_{\omega}^{p}(W)$ into $\mathcal{N}_{\omega \otimes \tau_{Z}}^{p+e}(W)\left(\right.$ resp. $\left.\mathcal{N}_{\omega \otimes \sigma_{Z}}^{p^{\prime}}(W)\right)$, where $p_{v}^{\prime}=\operatorname{Max}\left(0, p_{v}-e_{v}\right)$.

Proof. We first prove (3) in a formal sense when $|e|=1$. By (13.28) we have, for $u \in T$,

$$
\begin{align*}
&(\pi i)^{-1} D_{\omega, v}\left(q_{h}\left(\pi^{-1} i r\right) \mathbf{e}_{\mathbf{a}}^{n}(h z)\right)(u)=(\pi i)^{-1} \mathbf{e}_{\mathbf{a}}^{n}(h z)\left\{D_{v}\left(q_{h}\left(\pi^{-1} i r\right)\right)(u)\right.  \tag{*}\\
&\left.+2 \pi i \cdot \operatorname{tr}\left(h_{v} u_{v}\right) q_{h}\left(\pi^{-1} i r\right)+P_{\omega}^{v}(r, u) q_{h}\left(\pi^{-1} i r\right)\right\}
\end{align*}
$$

By (13.6c) we have

$$
(\pi i)^{-1} D_{\omega, v}\left(q_{h}\left(\pi^{-1} i r\right)\right)(u)=\sum_{\mu}\left(\partial q_{h} / \partial z_{v \mu}\right)\left(\pi^{-1} i r\right) \cdot\left(\pi^{-1} i r \cdot{ }^{t} u \cdot \pi^{-1} i r\right)_{v \mu}
$$

where $u_{v}=\sum_{\mu} u_{v \mu} a_{\mu}$ for $u \in T$ with a Q-rational C-basis $\left\{a_{\mu}\right\}$ of $T_{v}$, and in particular $z_{v}=\sum_{\mu} z_{v \mu} a_{\mu}$. Therefore the right-hand side of (*) can be written $\mathbf{e}_{\mathbf{a}}^{n}(h z) \ell_{h}\left(\pi^{-1} i r\right)(u)$ with a polynomial $\ell_{h}$ of degree $\leq p+v$ with values in $S_{1}\left(T_{v}, X\right)$. Write $\omega$ in the form $\omega(x)=\otimes_{v \in \mathbf{a}} \omega_{v}\left(x_{v \rho}, x_{v}\right)$, where $x_{v \rho}=x_{v}$ in Case SP. From (13.27) we easily see that $P_{\omega^{\varepsilon}}^{v^{\prime}}(r, u)=P_{\omega}^{v}\left(r^{[\varepsilon]}, u^{[\varepsilon]}\right)$ if $v \varepsilon=v^{\prime}$ on $F$ (even if $v \varepsilon=v^{\prime} \rho$ on $K$ ). Thus if we replace $\left(\omega, v, q_{h}\left(\pi^{-1} i r\right)\right)$ by $\left(\omega^{\varepsilon}, v^{\prime}, q_{h}^{\varepsilon}\left(\pi^{-1} i r^{[\varepsilon]}\right)\right)$, then $\ell_{h}\left(\pi^{-1} i r\right)(u)$ is replaced by $\ell_{h}^{\varepsilon}\left(\pi^{-1} i r^{[\varepsilon]}\right)\left(u^{[\varepsilon]}\right)$, as can easily be verified. This proves (3) for $D_{\omega, v}$. The general case of (3) can be proved by induction on $|e|$, which is not completely trivial. First, for $h \in S_{e}(T, X)$ put $g(u)=h\left(u^{[\varepsilon]}\right)$. Then $g \in S_{e \varepsilon}(T, X)$ and

$$
\begin{equation*}
\left[\left(\omega^{\varepsilon} \otimes \tau^{e \varepsilon}\right)(\alpha) g\right](u)=\left[\left(\omega \otimes \tau^{e}\right)^{\varepsilon}(\alpha) h\right]\left(u^{[\varepsilon]}\right) \text { for every } \alpha \in \mathfrak{K} . \tag{14.14}
\end{equation*}
$$

Here we can also take $\sigma$ in place of $\tau$. Notice that $\left(\omega^{\varepsilon} \otimes \tau^{e \varepsilon}\right)(\mathfrak{K})$ acts on $S_{e \varepsilon}(T, X)$, but $\left(\omega \otimes \tau^{e}\right)^{\varepsilon}(\mathfrak{K})$ acts on $S_{e}(T, X)$ (see $\left.\S 9.10\right)$. Now take $h \in C^{\infty}\left(\mathcal{H}, S_{e}(T, X)\right)$ and define $g \in C^{\infty}\left(\mathcal{H}, S_{e \varepsilon}(T, X)\right)$ by $g(u)=h\left(u^{[\varepsilon]}\right)$. Then from (14.14) we can easily derive

$$
\begin{equation*}
\left(D_{\omega^{\varepsilon} \otimes \tau^{e \varepsilon}}^{a} g\right)(u, w)=\left(D_{\left(\omega \otimes \tau^{e}\right)^{\varepsilon}}^{a} h\right)\left(u^{[\varepsilon]}, w\right) \quad(u, w \in T) \tag{14.15}
\end{equation*}
$$

Once this is established, our induction can be done in a straightforward way. This proves (3) for $D_{\omega}^{e}$. The statement concerning $E^{e}$ can be proved by using (13.10) in a similar and simpler way. Clearly our argument proves (4) for $Z=S_{e}(T)$. For $Z$ of a general type, we only have to observe that $\varphi_{Z}$ is $\mathbf{Q}$-rational, and so it sends $\mathcal{N}_{\omega \otimes \tau^{e}}^{e}(W)\left(\right.$ resp. $\left.\mathcal{N}_{\omega \otimes \sigma^{e}}^{e}(W)\right)$ into $\mathcal{N}_{\omega \otimes \tau_{Z}}^{e}(W)$ (resp. $\mathcal{N}_{\omega \otimes \sigma_{Z}}^{e}(W)$ ).

To prove (1) and (2), we use the symbols $\kappa, \zeta, B$, and $q$ in the proof of Theorem 14.9. By Lemma 6.17 we may assume that $q \in \mathcal{M}_{\kappa \mathbf{b}}(\mathbf{Q})$. Let $f \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$. In that proof we showed that $q f=\sum_{s \leq p} \pi^{-|s|} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{s}$ with $g_{s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\overline{\mathbf{Q}})$. Then for
every $\sigma \in \operatorname{Aut}(\mathbf{C} / \overline{\mathbf{Q}})$ we have $g_{s}^{\sigma}=g_{s}$, and hence $q f^{\sigma}=(q f)^{\sigma}=q f$ by (3). Thus $f^{\sigma}=f$, which means that $f$ is a $\overline{\mathbf{Q}}$-rational element of $\mathcal{N}_{\omega}^{p}$. Conversely let $f$ be a $\overline{\mathbf{Q}}$-rational element of $\mathcal{N}_{\omega}^{p}$. By Proposition 14.10 we have $f=\sum_{c \in B} c f_{c}$ with $f_{c} \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$. Let $q_{h}$ and $q_{c h}$ be the polynomials defined for $f$ and $f_{c}$ by (14.11). Then $q_{h}=\sum_{c \in B} c q_{c h}$. We have shown that $f_{c}$ is $\overline{\mathbf{Q}}$-rational, so that $q_{c h}$ is $\overline{\mathbf{Q}}$ rational. Then $q_{h}=q_{1 h}$, and hence $f=f_{1} \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$. This proves (1). To prove (2), we employ the expression $q f=\sum_{s \leq p}(\pi i)^{-|s|} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{s}$ with $g_{s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}$. Since (14.15) holds with $\sigma$ in place of $\tau$, we obtain, by (3),

$$
\begin{aligned}
{\left[(\pi i)^{-|s|} D_{\zeta \otimes \sigma^{s}}^{s} g_{s}\right]^{\varepsilon}\left(u^{[\varepsilon]}, w^{[\varepsilon]}\right) } & =\left[(\pi i)^{-|s|} D_{\left(\xi \otimes \sigma^{s}\right)^{\varepsilon}}^{s \varepsilon} g_{s}^{\varepsilon}\right]\left(u^{[\varepsilon]}, w\right) \\
& =\left[(\pi i)^{-|s|} D_{\zeta^{\varepsilon}}^{s \varepsilon} \otimes \sigma^{s \varepsilon} \ell_{s}\right](u, w),
\end{aligned}
$$

where $\ell_{s}$ is defined by $\ell_{s}(u)=g_{s}^{\varepsilon}\left(u^{[\varepsilon]}\right)$, and hence

$$
\begin{equation*}
\left[(\pi i)^{-|s|} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{s}\right]^{\varepsilon}=(\pi i)^{-|s|} \theta_{X}^{s \varepsilon} D_{\zeta^{\varepsilon} \otimes \sigma^{s \varepsilon}}^{s \varepsilon} \ell_{s} \tag{14.16}
\end{equation*}
$$

This shows that $(q f)^{\varepsilon}$ is meaningful as an element of $\mathcal{N}_{\varsigma^{\varepsilon}}^{p \varepsilon}$. If we take another $q^{\prime} \in \mathcal{M}_{\kappa \mathbf{b}}(\mathbf{Q})$, then $q\left(q^{\prime} f\right)^{\varepsilon}=\left(q q^{\prime} f\right)^{\varepsilon}=q^{\prime}(q f)^{\varepsilon}$, so that $(q f)^{\varepsilon} / q=\left(q^{\prime} f\right)^{\varepsilon} / q^{\prime}$. Define a function $f^{*}$ by $f^{*}=(q f)^{\varepsilon} / q$. This is defined where $q$ is not zero. Given $z_{0} \in \mathcal{H}$, we can find $q^{\prime} \in \mathcal{M}_{\kappa \mathbf{b}}(\mathbf{Q})$ so that $q^{\prime}\left(z_{0}\right) \neq 0$. Since $f^{*}=\left(q^{\prime} f\right)^{\varepsilon} / q^{\prime}$, we see that $f^{*}$ is defined as a $C^{\infty}$ function on $\mathcal{H}$. Besides, the equality $f^{*}=(q f)^{\varepsilon} / q$ shows that $f^{*} \in \mathcal{N}_{\omega^{\varepsilon}}^{p \varepsilon}$, and $q f^{*}=(q f)^{\varepsilon}$. Now $(q f)^{\varepsilon}$ and $q f^{\varepsilon}$ coincide as formal series, and hence $q f^{*}$ coincides with $q f^{\varepsilon}$ as a formal series. As shown in $\S 5.10, q$ is not a zero-divisor in the ring of formal series defined there. Therefore, expressing $f^{\varepsilon}$ and $f^{*}$ as polynomials in $r$ whose coefficients are elements of the ring, we see that $f^{\varepsilon}$ is the expansion of $f^{*}$ in the sense of (14.11). This proves (2).
14.13. Proposition (Cases SP and UT). Let $\Phi$ be the Galois closure of $K$ over Q. Then the following assertions hold:
(1) $\mathcal{N}_{\omega}^{p}=\mathcal{N}_{\omega}^{p}(\Phi) \otimes_{\Phi} \mathbf{C}$ and $\mathcal{N}_{\omega}^{p}(\Gamma)=\mathcal{N}_{\omega}^{p}(\Gamma, \Phi) \otimes_{\Phi} \mathbf{C}$ for every $\Gamma \subset G$ as in Theorem 10.4 such that $\mathcal{M}_{\mu \mathbf{b}}(\Gamma, \overline{\mathbf{Q}}) \neq\{0\}$ for some $\mu \in \mathbf{Z},>0$.
(2) If $W$ is a subfield of $\overline{\mathbf{Q}}$ containing $\Phi$ and $\mathbf{Q}_{\mathrm{ab}}$, then $\mathcal{N}_{\omega}^{p}(W)$ is stable under the map $f \mapsto f \|_{\omega} \alpha$ for every $\alpha \in G$.

Proof. Take an arbitrary subfield $\Psi$ of $\overline{\mathbf{Q}}$ containing $\Phi$; let $B$ be a $\Psi$-basis of $\mathbf{C}$ and $M$ a Q-rational C-basis of $\bigcup_{e \leq p} S_{e}(T, X)$. Suppose $0=\sum_{c \in A} c f_{c}$ with a finite subset $A$ of $B$ and $f_{c} \in \mathcal{N}_{\omega}^{p}(\Psi)$. Put $f_{c}=\sum_{\ell \in M} \sum_{h \in S} a(c, \ell, h) \ell\left(\pi^{-1} i\right.$. $r(z)) \mathbf{e}_{\mathbf{a}}^{n}(h z)$ with $a(c, \ell, h) \in \Psi$. Then $0=\sum_{c \in A} c a(c, \ell, h)$ for every $(\ell, h)$, and so $a(c, \ell, h)=0$ for every $(c, \ell, h)$, that is, $f_{c}=0$ for every $c$. This result is of course applicable to the case $\Psi=\Phi$. Thus, to prove (1), it is sufficient to prove, in view of Proposition 14.10, that $\mathcal{N}_{\omega}^{p}(\Phi)$ spans $\mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$ over $\overline{\mathbf{Q}}$. Let $f \in \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$. We use the expression $q f=\sum_{s \leq p}(\pi i)^{-|s|} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{s}$ with $g_{s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\overline{\mathbf{Q}})$ already employed in the proof of Theorem 14.12. By Theorem 10.4 (4) we have $g_{s}=\sum_{c \in A} c g_{c, s}$ with a finite subset $A$ of $\overline{\mathbf{Q}}$, linearly independent over $\Phi$, and with $g_{c, s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\Phi)$. Then $q f=\sum_{c \in A} c k_{c}$ with $k_{c}=\sum_{s \leq p}(\pi i)^{-|s|} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{c, s}$. By Theorem 14.12 (4), $k_{c} \in \mathcal{N}_{\zeta}^{p}(\Phi)$. If we change $q$ for another $q^{\prime}$, then we have a similar expression $q^{\prime} f=\sum_{c \in A} c k_{c}^{\prime}$. Then $0=\sum_{c \in A} c\left(q^{\prime} k_{c}-q k_{c}^{\prime}\right)$, and hence $q^{\prime} k_{c}=q k_{c}^{\prime}$. Thus $q^{-1} k_{c}=q^{\prime-1} k_{c}^{\prime}$. Since we can choose $q \in \mathcal{M}_{\kappa \mathbf{a}}(\mathbf{Q})$ so that $q$ does not vanish at any given point of $\mathcal{H}$, we can define a $C^{\infty}$ function $f_{c}$ on $\mathcal{H}$ by $f_{c}=q^{-1} k_{c}$. Clearly $f_{c} \in \mathcal{N}_{\omega}^{p}(\Phi)$, and $f=\sum_{c \in A} c f_{c}$. This proves the first part of (1). Suppose
$f \in \mathcal{N}_{\omega}^{p}(\Gamma, \overline{\mathbf{Q}})$. Taking $q$ to be a power of an element of $\mathcal{M}_{\mu \mathbf{b}}(\Gamma, \overline{\mathbf{Q}})$, we may assume, by (14.9a), that $g_{s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\Gamma, \overline{\mathbf{Q}})$; thus, by Theorem 10.4 (4) we may assume that $g_{c, s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\Gamma, \Phi)$. Then $k_{c} \in \mathcal{N}_{\zeta}^{p}(\Gamma, \Phi)$, and hence $f_{c} \in \mathcal{N}_{\omega}^{p}(\Gamma, \Phi)$, which, together with the second part of Proposition 14.10 proves the second part of (1).

To prove (2), we again take $f \in \mathcal{N}_{\omega}^{p}(W)$ and apply the above procedure to $f$. We can take $A$ to be the set $\left\{a b \mid a \in A_{1}, b \in A_{2}\right\}$, where $A_{1}$ is a finite subset of $W$ including 1 and linearly independent over $\Phi$, and $A_{2}$ is a finite subset of $\overline{\mathbf{Q}}$ including 1 and linearly independent over $W$. Then, as shown above, $f=\sum_{a \in A_{1}} \sum_{b \in A_{2}} a b f_{a b}$ with $f_{a b} \in \mathcal{N}_{\omega}^{p}(\Phi)$. Since $f \in \mathcal{N}_{\omega}^{p}(W)$, what we said at the beginning shows that $f_{a b}=0$ if $b \neq 1$, and so $f=\sum_{a \in A_{1}} a f_{a}$. Now we obtained $f_{a}$ in the form $f_{a}=q^{-1} \sum_{s \leq p}(\pi i)^{-|s|} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{a, s}$ with $g_{a, s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}(\Phi)$. Apply $\|_{\omega} \alpha$ to this equality with $\alpha \in G$. Then from Theorem 9.13 (3), (12.21), (14.2), and Theorem 14.12 (4) we see that $f_{a} \|_{\omega} \alpha \in \mathcal{N}_{\omega}^{p}\left(\Phi \mathbf{Q}_{\mathrm{ab}}\right)$. This completes the proof.
14.14. Let us now extend our results of Sections 12 through 14 to the case of modular forms of half-integral weight. We naturally restrict our discussion to Type CT (that is, Case SP), employing the symbols introduced in $\S 10.6$; in particular we take a symbol $\psi$, which we called a quasi-representation of $\mathfrak{K}$ in $\S 10.6$, given by

$$
\begin{equation*}
\psi(x)=\operatorname{det}(x)^{\mathbf{a} / 2} \omega(x) \tag{14.17}
\end{equation*}
$$

with a Q-rational representation $\{\omega, X\}$ of $G L_{n}(\mathbf{C})^{\mathbf{a}}$. (If $\omega(x)=\operatorname{det}(x)^{m}$ with $m \in \mathbf{Z}^{\mathbf{a}}$, then we can put $\psi(x)=\operatorname{det}(x)^{k}$ with $k_{v}=m_{v}+(1 / 2)$; see $\S 6.10$.) We also consider the group $\mathcal{G}$ consisting of all couples $(\alpha, p)$ formed by $\alpha \in G$ and a holomorphic function $p$ on $\mathcal{H}$ such that $p(z)^{2} / j_{\alpha}^{\mathrm{a}}$ is a root of unity, the grouplaw being defined by $(\alpha, p)\left(\alpha^{\prime}, p^{\prime}\right)=\left(\alpha \alpha^{\prime}, p\left(\alpha^{\prime} z\right) p^{\prime}(z)\right)$. For $\gamma$ in the group $\Gamma^{\theta}$ of (6.30) we put $\tilde{\gamma}=\left(\gamma, h_{\gamma}\right)$ with $h_{\gamma}$ of Theorem 6.8. Then $\gamma \mapsto \widetilde{\gamma}$ is an injective homomorphism of $\Gamma^{\theta}$ into $\mathcal{G}$. We call a subgroup $\Delta$ of $\mathcal{G}$ a congruence subgroup of $\mathcal{G}$ if the projection map of $\mathcal{G}$ onto $G$ gives an isomorphism of $\Delta$ onto a congruence subgroup $\Gamma$ of $G$, and the inverse of this isomorphism coincides with the map $\gamma \mapsto \widetilde{\gamma}$ on a congruence subgroup of $\Gamma^{\theta} \cap \Gamma$. Any conjugate $\xi \Delta \xi^{-1}$ of such a $\Delta$ with $\xi \in \mathcal{G}$ is also a congruence subgroup, by virtue of Theorem 6.9 (1). We shall often view a congruence subgroup $\Gamma$ of $\Gamma^{\theta}$ as a congruence subgroup of $\mathcal{G}$ by identifying it with its image in $\mathcal{G}$ through the map $\gamma \mapsto \widetilde{\gamma}$.

Now for $\alpha=\left(\alpha_{0}, p\right) \in \mathcal{G}$ and $f \in \mathbf{C}^{\infty}(\mathcal{H}, X)$ we put

$$
\begin{equation*}
\left(f \|_{\psi} \alpha\right)(z)=p(z)^{-1}\left(f \|_{\omega} \alpha_{0}\right)(z) \tag{14.18a}
\end{equation*}
$$

If $\omega(x)=\operatorname{det}(x)^{m}$ and $\psi(x)=\operatorname{det}(x)^{k}$ with $k=m+\mathbf{a} / 2$ as above, then we write $f \|_{k} \alpha$ for $f \|_{\psi} \alpha$; then

$$
\begin{equation*}
\left(f \|_{k} \alpha\right)(z)=p(z)^{-1}\left(f \|_{m} \alpha_{0}\right)(z) \tag{14.18b}
\end{equation*}
$$

For simplicity we put $M_{\alpha}(z)=M_{\alpha_{0}}(z)$ and $\alpha z=\alpha_{0} z$. We also define quasirepresentations $\psi \otimes \tau^{e}$ and $\psi \otimes \sigma^{e}$ by

$$
\begin{equation*}
\left(\psi \otimes \tau^{e}\right)(x)=\operatorname{det}(x)^{\mathbf{a} / 2}\left(\omega \otimes \tau^{e}\right)(x), \quad\left(\psi \otimes \sigma^{e}\right)(x)=\operatorname{det}(x)^{\mathbf{a} / 2}\left(\omega \otimes \sigma^{e}\right)(x) \tag{14.19}
\end{equation*}
$$

We define $D_{\psi, v}^{e}, D_{\psi, v}, D_{\psi}^{e}$, and $D_{\psi}^{Z}$ by (12.17), (12.18), (13.21) and (13.22) with $\psi$ in place of $\rho$, in which we take $\psi(\Xi)=\operatorname{det}(\eta)^{\mathbf{a} / 2} \omega(\Xi)$ with positive $\operatorname{det}(\eta)^{\mathbf{a} / 2}$. Then Proposition 12.10 and (12.21) are true with $\psi$ in place of $\rho$ and with $\alpha \in \mathcal{G}$. The proof of Proposition 12.10 is valid for $\psi$ if we put

$$
\begin{equation*}
\psi\left(M_{\alpha}\right)=p \cdot \omega\left(M_{\alpha}\right) \tag{14.20}
\end{equation*}
$$

Lemma 13.9 is valid also for half-integral $k$ and $\alpha \in \mathcal{G}$.
We already defined $\mathcal{M}_{\psi}, \mathcal{S}_{\psi}$, and $\mathcal{A}_{\psi}$ in $\S 10.6$. We denote by $C_{\psi}(\Delta)$ the set of all $f \in C^{\infty}(\mathcal{H}, X)$ such that $f \|_{\psi} \gamma=f$ for every $\gamma \in \Delta$, and by $C_{\psi}$ the union of $C_{\psi}(\Delta)$ for all congruence subgroups $\Delta$ of $\mathcal{G}$. We denote by $\mathcal{N}_{\psi}^{p}(\Delta)$ the subset of $\mathcal{N}^{p}(\mathcal{H}, X) \cap C_{\psi}(\Delta)$ consisting of the functions satisfying the cusp condition, which is required only when $G$ is isogenous to $S L_{2}(\mathbf{Q})$, and which is an obvious modification of (13.18a). We then denote by $\mathcal{N}_{\psi}^{p}$ the union of $\mathcal{N}_{\psi}^{p}(\Delta)$ for all such $\Delta$. The inner product $\langle f, g\rangle$ for $f, g \in C_{\psi}$ can be defined by (12.35a) with $\psi$ in place of $\rho$. Then Theorem 12.15 is valid with $\psi$ in place of $\rho$, since the problem can be formulated on a suitable covering group of $G_{\mathbf{a}}$, and can be reduced to the results of [S90] by virtue of the principle of [S94b, Proposition 2.2].

Now $\psi$ can be viewed as a local homomorphism of $\mathfrak{K}$ into $G L(X)$, and so $d \psi$ : $\left(\mathbf{C}_{n}^{n}\right)^{\mathbf{a}} \rightarrow \operatorname{End}(X)$ is meaningful. Therefore Lemma 13.14 and Proposition 13.15 are valid with $\psi$ in place of $\rho$.

We say that an element $f$ of $\mathcal{N}_{\psi}^{p}$ is arithmetic at $w \in \mathcal{H}_{\mathrm{CM}}$ if $\mathfrak{P}_{\psi}(w)^{-1} f(w)$ is $\overline{\mathbf{Q}}$-rational, where we naturally put $\mathfrak{P}_{\psi}(w)=\prod_{v \in \mathbf{a}} \operatorname{det}\left(\mathfrak{p}_{v}(w)\right)^{1 / 2} \mathfrak{P}_{\omega}(w)$. Since we are interested only in the algebraicity, the choice of square roots does not matter. We then dfine $\mathcal{N}_{\psi}^{p}(\overline{\mathbf{Q}})$ to be the set of all $f \in \mathcal{N}_{\psi}^{p}$ that are arithmetic at every $w \in \mathcal{H}_{\mathrm{CM}}$. Now for $f \in \mathcal{N}_{\psi}^{p}$ we have an expansion of type (14.10), and so we define $f^{\varepsilon}$ by (14.13) and the rationality of $f$ over a subfield of $\mathbf{C}$ in the same way as for the elements of $\mathcal{N}_{\omega}^{p}$. Then we can verify that the results up to Proposition 14.13 are all valid with $\psi$ in place of $\omega$.
14.15. Let us now specialize our discussion to the Hilbert modular case, by taking $G=S L_{2}(F)$ and $X=\mathbf{C}$. Then we can put $T=\mathbf{C}^{\mathbf{a}}, S_{e}(T, X)=\mathbf{C}$ by identifying $h \in S_{e}(T, X)$ with $h(1, \ldots, 1)$, and $\omega(x)=x^{k}$ for $x \in\left(\mathbf{C}^{\times}\right)^{\mathbf{a}}$ with an integral or a half-integral weight $k$ (see $\S \S 6.10$ and 14.13). Therefore, rewriting the expansion of (14.11), we see that every element of $\mathcal{N}_{k}^{p}$ is of the form

$$
\begin{equation*}
f(z)=\sum_{h \in F} \sum_{0 \leq e \leq p} c(h, e)(\pi y)^{-e} \mathbf{e}_{\mathbf{a}}(h z) \quad\left(z \in \mathfrak{H}_{1}^{\mathbf{a}}\right) \tag{14.21}
\end{equation*}
$$

with $c(h, e) \in \mathbf{C}$, where $y=\operatorname{Im}(z) ; c(h, e) \neq 0$ only if $h=0$ or $h \gg 0$. Then

$$
\begin{equation*}
f^{\varepsilon}(z)=\sum_{h \in F} \sum_{0 \leq e \leq p} c(h, e)^{\varepsilon}(\pi y)^{-e \varepsilon} \mathbf{e}_{\mathbf{a}}(h z) . \tag{14.22}
\end{equation*}
$$

Thus $f \in \mathcal{N}_{k}^{p}(W)$ with a subfield $W$ of $\mathbf{C}$ if and only if $c(h, e) \in W$ for every $h$ and $e$. Write $D_{k}^{e}$ for $D_{\omega}^{e}$. Then $(\pi i)^{-|e|} D_{k}^{e}=(\pi i)^{-|e|} \prod_{v \in \mathbf{a}} \delta_{k_{v}}^{e_{v}}$ with the symbol $\delta$ of (12.39). Also, we can ignore ( $u$ ) and ( $u^{[\varepsilon]}$ ) in Theorem 14.12 (3).

We note that the space $\mathcal{N}_{k}^{p}$ can be completely determined; see [ S 87 , Theorem 5.2]. Also a generalization of Theorem 10.9 (1) can be given as follows:
14.16. Theorem. Let $Y, h, w$, and $Y^{*}$ be as in §9.4; suppose $G=S L_{2}(F)$. (Thus $Y$ is a $C M$-field and $h$ is an $F$-linear ring-injection of $Y$ into $F_{2}^{2}$.) Let $f \in$ $\mathcal{M}_{k}(\overline{\mathbf{Q}})$ and $p \in \mathcal{M}_{k+2 e}(\overline{\mathbf{Q}})$ with $k \in \mathbf{Z}^{\mathbf{a}}$ and $0 \leq e \in \mathbf{Z}^{\mathbf{a}}$. Given $\sigma \in \operatorname{Gal}(\overline{\mathbf{Q}} / \mathbf{Q})$, let $b$ be an element of $\left(Y^{*}\right)_{\mathbf{A}}$ such that $\sigma=\left[b, Y^{*}\right]$ on $Y_{\mathrm{ab}}^{*}$ and let $r=h\left(g(b)^{-1}\right)$ with $g$ of (9.3). Then

$$
\left\{p^{-1}(\pi i)^{-|e|} D_{k}^{e} f\right\}(w)^{\sigma}=p^{(r, \sigma)}(w)^{-1}(\pi i)^{-|e|}\left(D_{k}^{e} f^{(r, \sigma)}\right)(w)
$$

for every $C M$-point $w$ such that $p(w) \neq 0$.
For the proof, see [S75, Main Theorem III]. Notice that $p^{(r, \sigma)}(w) \neq 0$ by Theorem 10.9 (3).

The arithmeticity of $D_{\omega}^{e} f$ at CM-points was first obtained in [S75] for $\widetilde{G}=$ $G L_{2}(F)$ and in [S80] for orthogonal and unitary groups. A general framework of the arithmeticity problems of zeta functions and Eisenstein series was presented in [S78c]. The notion of near holomorphy was introduced in [S86], and further developed in [S87a]. Most of the results in Sections 13 and 14 were essentially given in those two papers. A notable exception is Theorem 14.12, which was not given there. As explained in [S87a, Section 3], we can develop the theory axiomatically so that all the results for other types of groups can be proved in the same manner, once a few basic facts such as the result corresponding to Proposition 14.5 in each case is established. Such a result was given in [S78c, Theorem 4] for a certain quaternion unitary group belonging to Type C, and also in [S80, Proposition 6.6] in the orthogonal case. We can also characterize the elements of $\mathcal{M}_{\omega}(\overline{\mathbf{Q}})$ in Case UB by the properties of the theta functions that appear as the Fourier coefficients of a given automorphic form; see [S78a] and [S78b, §7]. If $\operatorname{dim}\left(H_{v}\right) \leq 1$ for every $v \in \mathbf{a}$ (which is the case if $G=S L_{2}(F)$ ), then we can completely determine the structure of $\mathcal{N}_{k}^{p}$, so that we can state Proposition 14.2 in a much stronger form; for details, see [S87a, Theorems 5.2 and 5.5].

## 15. Holomorphic projection

15.1. Our next aim is to find a certain projection map $\mathcal{N}_{\omega}^{p} \rightarrow \mathcal{M}_{\omega}$. It is necesssary to consider $\theta h=\sum_{\nu \in N} h\left(a_{\nu}, b_{\nu}\right)$ defined by (12.33) with any pair of dual bases $\left\{a_{\nu}\right\}$ and $\left\{b_{\nu}\right\}$ of $T$ for several specific $h \in M l_{2}(T, X)=S_{1}\left(T, S_{1}(T, X)\right)$. (We again fix one $v \in \mathbf{a}$, and drop the subscript $v$ from the objects $T_{v}, \tau_{v}^{p}$, etc.) For example, for $h(x, y)={ }^{t} x y$ we can easily verify that

$$
\begin{equation*}
\theta h=\sum_{\nu \in N}{ }^{t} a_{\nu} b_{\nu}=\lambda(T) 1_{n}, \tag{15.1}
\end{equation*}
$$

where $\lambda(T)=m$ for Type A and $\lambda(T)=(n+1) / 2$ for Type C.
We now define a C-linear endomorphism $\psi$ of $S_{p}(T)$ by

$$
\begin{gather*}
\psi=0 \quad \text { if } p=1  \tag{15.2a}\\
(\psi h)(x)=\sum_{\nu \in N} h_{*}\left(a_{\nu}, x^{\cdot t} b_{\nu} x, x, \ldots, x\right) \quad(p>1) \tag{15.2b}
\end{gather*}
$$

for $h \in S_{p}(T)$ and $x \in T$. We can easily verify that this is independent of the choice of dual bases, and $\psi \tau^{p}(\alpha)=\tau^{p}(\alpha) \psi$ for every $\alpha \in K^{c}$, and hence, for each irreducible subspace $Z$ of $S_{p}(T)$ there is a constant $c_{Z}$ such that $\psi h=c_{Z} h$ for every $h \in Z$. Thus $c_{Z}=0$ if $p=1$.
15.2. Lemma. The constant $c_{Z}$ is a rational number such that $-1 \leq c_{Z} \leq 1$ for Type $A$ and $-1 / 2 \leq c_{Z} \leq 1$ for Type C. Moreover, $c_{Z}=1$ if $Z$ contains the element $h$ of $S_{p}(T)$ given by $h(x)=x_{11}^{p}$, and $c_{Z}<1$ otherwise. In particular, $c_{Z}=-1$ (resp. $c_{Z}=-1 / 2$ ) for Type $A$ (resp. Type $C$ ) if $Z$ contains the element $h$ of $S_{2}(T)$ given by $h(x)=\operatorname{det}_{2}(x)$.

Proof. Type A. Let $\ell=\operatorname{Min}(m, n)$ and let $\zeta=\sum_{i=1}^{\ell} e_{i i}(\in T)$ with the standard matrix units $e_{i j}$. By Theorem 12.7 we can take a highest weight vector of
$Z$ in the form $h(x)=\prod_{i=1}^{\ell} \operatorname{det}_{i}(x)^{c_{i}}$ with $0 \leq c_{i} \in \mathbf{Z}$. Take $k \in M l_{p}(T, \mathbf{C})$ so that $h(x)=k(x, \ldots, x)$ by the rule which can be illustrated by the following example: if $p=9, c_{1}=c_{2}=2, c_{3}=1$, then

$$
k(r, s, t, u, v, w, x, y, z)=r_{11} s_{11}\left|\begin{array}{cc}
t_{11} & u_{12} \\
t_{21} & u_{22}
\end{array}\right| \cdot\left|\begin{array}{cc}
v_{11} & w_{12} \\
v_{21} & w_{22}
\end{array}\right| \cdot\left|\begin{array}{lll}
x_{11} & y_{12} & z_{13} \\
x_{21} & y_{22} & z_{23} \\
x_{31} & y_{32} & z_{33}
\end{array}\right|
$$

Then $p!h_{*}=\sum_{\pi} k_{\pi}$ with $k_{\pi}\left(x_{1}, \ldots, x_{p}\right)=k\left(x_{\pi(1)}, \ldots, x_{\pi(p)}\right)$. Since $h(\zeta)=1$, we have $c_{Z}=(\psi h)(\zeta)=(1 / p!) \sum_{\pi} \sum_{\nu} k_{\pi}\left(a_{\nu}, \zeta{ }^{t} b_{\nu} \zeta, \zeta, \ldots, \zeta\right)$. Each $k_{\pi}(x, y$, $\zeta, \ldots, \zeta)$ belongs to the following three types of functions: $x_{i i} y_{i i}, x_{i i} y_{j j}(i \neq j)$, and $x_{i i} y_{j j}-x_{j i} y_{i j}(i \neq j)$. The value $\sum_{\nu} k_{\pi}\left(a_{\nu}, \zeta \cdot{ }^{t} b_{\nu} \zeta ; \zeta, \ldots, \zeta\right)$ is 1,0 , and -1 , respectively. Therefore we obtain $-1 \leq c_{Z} \leq 1$. If $h(x)=x_{11}^{p}$, then the only possible type is $x_{11} y_{11}$, and hence $c_{Z}=1$. If $\operatorname{det}_{i}(x)$ with $i>1$ is involved in $h$, then $x_{11} y_{22}-x_{21} y_{12}$ can always occur, and hence $c_{Z}<1$. In particular, if $h(x)=\operatorname{det}_{2}(x)$, then that is the only possible type, so that $c_{Z}=-1$.

Type C. We can employ the same technique as for Type A with $1_{n}$ in place of $\zeta$, $h(x)=\prod_{i=1}^{n} \operatorname{det}_{i}(x)^{c_{i}}$ and the same $k$. Then $p!c_{Z}=\sum_{\pi} \sum_{\nu} k_{\pi}\left(a_{\nu}, b_{\nu}, 1, \ldots, 1\right)$. Each $k_{\pi}(x, y, 1, \ldots, 1)$ belongs to the following three types of functions: $x_{i i} y_{i i}$, $x_{i i} y_{j j}(i \neq j)$, and $x_{i i} y_{j j}-x_{i j} y_{i j}(i \neq j)$. The value $\sum_{\nu} k_{\pi}\left(a_{\nu}, b_{\nu}, 1, \ldots, 1\right)$ is 1,0 , and $-1 / 2$, respectively, and hence $-1 / 2 \leq c_{Z} \leq 1$. The remaining part concerning $c_{Z}$ for each $Z$ can be proved in the same manner as for Type A.

Before stating the next proposition, we define an operator $L_{\omega, v}$ by

$$
\begin{equation*}
L_{\omega, v}=-\theta D_{\omega \otimes \sigma_{v}^{1}, v} E_{v} \tag{15.3}
\end{equation*}
$$

This is a special case of (12.31), which we already mentioned in Corollary 12.16 and Lemma 13.16, and $L_{\omega, v}\left(f \|_{\omega} \alpha\right)=\left(L_{\omega, v} f\right) \|_{\omega} \alpha$ for every $\alpha \in G_{\mathbf{a}}$.
15.3. Proposition. Let $0 \neq p=\left(p_{v}\right)_{v \in \mathbf{a}} \in \mathbf{Z}^{\mathbf{a}}$ with $p_{v} \geq 0$ for every $v \in \mathbf{a}$ and $\{\omega, X\}$ a representation or a quasi-representation (in the sense of §14.13) of $\mathfrak{K}$. Suppose that for every $v \in \mathbf{a}$ such that $H_{v}$ is nontrivial we have $\omega\left(a_{v}, b_{v}\right)=$ $\operatorname{det}\left(b_{v}\right)^{k_{v}}$ with $k_{v} \in(1 / 2) \mathbf{Z}$. For an irreducible subspace $Z$ of $S_{i}\left(T_{v}\right)$, put

$$
\alpha_{Z}=i\left\{k_{v}-\kappa+(1-i) c_{Z}\right\}
$$

where $\kappa=m+n$ or $n+1$ according as $G$ is of Type $A$ or $C$, where we understand that $\kappa=2 n$ for Type AT. Suppose that for each $v$ such that $p_{v}>0$ the number $k_{v}$ satisfies the following inequalities:

$$
\begin{gathered}
k_{v}>m+n+p_{v}-1 \quad \text { or } \quad k_{v}<m+n+1-p_{v} \quad \text { if } G \text { is of Type } A, \\
k_{v}>n+p_{v} \quad \text { or } \quad k_{v}<n+\left(3-p_{v}\right) / 2 \quad \text { if } G \text { is of Type } C .
\end{gathered}
$$

Put $A_{v}^{i}=\prod_{Z}\left(1-\alpha_{Z}^{-1} L_{\omega, v}\right)$ for $0<i \leq p_{v}$, where $Z$ runs over all the irreducible subspaces of $S_{i}\left(T_{v}\right)$, and $\mathfrak{A}=\prod_{v \in \mathbf{a}^{*}} \prod_{i=1}^{p_{v}} A_{v}^{i}$, where $\mathbf{a}^{*}=\left\{v \in \mathbf{a} \mid p_{v}>0\right\}$. (Notice that the estimate of $c_{Z}$ given in Lemma 15.2 shows that $\alpha_{Z} \neq 0$.) Let $f \in \mathcal{N}_{\omega}^{p}$. Then $\mathfrak{A} f \in \mathcal{M}_{\omega}$ and $f=\mathfrak{A} f+\sum_{v \in \mathbf{a}^{*}} L_{\omega, v} t_{v}$ with $t_{v} \in \mathcal{N}_{\omega}^{p}$.

Proof. We fix one $v$ and consider $f$ as a function of $z_{v}$, suppressing the remaining variables. By (13.16) we have $f=\sum_{i=0}^{q} f_{i}, f_{i}\left(z_{v}\right)=g_{i}\left(r_{v}\left(z_{v}\right)\right), q=p_{v}$, with a holomorphic map $g_{i}: H_{v} \rightarrow S_{i}\left(T_{v}\right)$. Then $E_{v} f=\sum_{i=1}^{q} E_{v} f_{i}$, and $E_{v} f_{i}$ is of degree $i-1$ in $r_{v}$. To study the highest term $E_{v} f_{q}$, write simply $g$ for $g_{q}$. Then $g_{*}$ is a holomorphic map of $H_{v}$ into $M l_{q}\left(T_{v}, X\right)$. Let us now write simply $T, r, k, D, E$, etc. for $T_{v}, r_{v}, k_{v}, D_{v}, E_{v}$, etc. By (13.10), for $u \in T$ we have

$$
\left(E f_{q}\right)(u)=-\sum_{\nu \in N} u_{\nu}\left(\partial / \partial r_{\nu}\right) g_{*}(r, \ldots, r)=-q g_{*}(u, r, \ldots, r)
$$

By (12.7b) and (12.18) we see that $\left(D_{\omega \otimes \sigma^{1}} E f_{q}\right)(u, v)=\operatorname{det}(\eta)^{-k} Y\left(\xi u \cdot{ }^{t} \eta, v\right)$ with $Y(u, v)=D\left\{\operatorname{det}(\eta)^{k}\left(E f_{q}\right)\left(\xi^{-1} u \cdot{ }^{t} \eta^{-1}\right)\right\}(v)$ for $u, v \in T$. Then by (13.8a) and (13.6a, b) we have $D\left(\operatorname{det}(\eta)^{k}\right)(v)=k \cdot \operatorname{det}(\eta)^{k} \operatorname{tr}\left({ }^{t} r v\right)$, and

$$
D\left(\xi^{-1} u \cdot{ }^{t} \eta^{-1}\right)(v)=-r \cdot{ }^{t} v \xi^{-1} u \cdot{ }^{t} \eta^{-1}-\xi^{-1} u \cdot{ }^{t} \eta^{-1} \cdot t v r
$$

Therefore

$$
\begin{aligned}
& -q^{-1}\left(D_{\omega \otimes \sigma^{1}} E f_{q}\right)(u, v)=k \cdot \operatorname{tr}\left({ }^{t} r v\right) g_{*}(u, r, \ldots, r)-g_{*}(w, r, \ldots, r) \\
& \quad+(q-1) g_{*}(u,(D r)(v), r, \ldots, r)+\sum_{\nu \in N} v_{\nu}\left(\partial g_{*} / \partial z_{\nu}\right)(u, r, \ldots, r)
\end{aligned}
$$

where $w=r \cdot{ }^{t} v u+u \cdot{ }^{t} v r$. Applying $\theta$ to this equality, we obtain $k g(r)$ from the first term on the right-hand side and, by (15.1), $-\kappa g(r)$ from the second term with $\kappa$ given as in our proposition. Now $(D r)(v)$ is given by (13.6c), and therefore $\theta$ times the third term is $(1-q)(\psi g)(r)$. The last sum $\sum_{\nu}$ is of degree at most $q-1$ in $r_{v}$. Thus we obtain

$$
L_{\omega, v} f \equiv L_{\omega, v} f_{q} \equiv p_{v}\left\{k_{v}-\kappa+\left(1-p_{v}\right) \psi\right\} g(r) \quad\left(\bmod \mathcal{N}^{p^{\prime}}\right)
$$

where $p_{v}^{\prime}=p_{v}-1$ and $p_{t}^{\prime}=p_{t}$ for $v \neq t \in \mathbf{a}$. (This is true even if $p_{v}=1$.) Let $\varphi_{Z}$ be the projection map $S_{q}\left(T_{v}\right) \rightarrow Z$. Then $p_{v}\left\{k_{v}-\kappa+\left(1-p_{v}\right) \psi\right\}=\sum_{Z} \alpha_{Z} \varphi_{Z}$. Now $L_{\omega, v}$ maps $\mathcal{N}_{\omega}^{p}$ into itself. Therefore, for an irreducible subspace $W$ of $S_{q}\left(T_{v}\right)$ we have

$$
\left(1-\alpha_{W}^{-1} L_{\omega, v}\right) f \equiv \sum_{Z}\left(1-\alpha_{W}^{-1} \alpha_{Z}\right) \varphi_{Z} g \equiv \sum_{Z \neq W}\left(a_{Z} \varphi_{Z} g\right)(r) \quad\left(\bmod \mathcal{N}^{p^{\prime}}\right)
$$

with $a_{Z}=1-\alpha_{W}^{-1} \alpha_{Z}$. Call the last sum $g^{\prime}$. Taking $\left(1-\alpha_{W}^{-1} L_{\omega, v}\right) f$ and $g^{\prime}$ in place of $f$ and $g$, we obtain, for another irreducible subspae $Y$ of $S_{q}\left(T_{v}\right)$,

$$
\left(1-\alpha_{Y}^{-1} L_{\omega, v}\right)\left(1-\alpha_{W}^{-1} L_{\omega, v}\right) f \equiv \sum_{Z \notin\{Y, W\}}\left(b_{Z} \varphi_{Z} g^{\prime}\right)(r) \quad\left(\bmod \mathcal{N}^{p^{\prime}}\right)
$$

with $b_{Z} \in \mathbf{Q}$. Repeating this procedure, we find that $A_{v}^{q} f \in \mathcal{N}_{\omega}^{p^{\prime}}$ with $A_{v}^{q}$ defined in our proposition. Therefore if $h=\mathfrak{A} f$, then $h \in \mathcal{N}_{\omega}^{0}=\mathcal{M}_{\omega}$. Since $\prod_{i} A_{v}^{i}$ is a polynomial in $L_{\omega, v}$ whose constant term is 1 , we obtain the desired expression for $f$.
15.4. Corollary. Let $f \in \mathcal{N}_{\omega}^{p}$ and $h=\mathfrak{A} f$ with $\mathfrak{A}$ as in Proposition 15.3. Then we have $\langle\varphi, f\rangle=\langle\varphi, h\rangle$ for $\varphi \in \mathcal{M}_{\omega}$, provided either $\varphi$ is a cusp form, or $\Gamma \backslash \mathcal{H}$ is compact for a congruence subgroup $\Gamma$.

Proof. We have $f=h+\sum_{v \in \mathbf{a}^{\prime}} L_{\omega, v} t_{v}$ with $t_{v} \in \mathcal{N}_{\omega}^{p}$. By Theorem 12.15,

$$
\left\langle\varphi, L_{\omega, v} t_{v}\right\rangle=\left\langle\varphi,-\theta D_{\omega \otimes \sigma_{v}^{1}, v} E_{v} t_{v}\right\rangle=\left\langle E_{v} \varphi, E_{v} t_{v}\right\rangle=0
$$

since $E_{v} \varphi=0$ because of the holomorphy of $\varphi$. This proves the desired equality. We need a suitable convergence condition, which is certainly satisfied for $\varphi$ or $\Gamma$ as above.
15.5. In Cases SP and UT, for $f \in \mathcal{N}_{\omega}^{p}$ let $q_{h}(u, f)$ denote the polynomial in $u \in T$ which appears as a Fourier coefficient of $f$ as in (14.11). Then we put

$$
\begin{gather*}
\mathcal{R}_{\omega}^{p}=\left\{f \in \mathcal{N}_{\omega}^{p} \mid q_{h}(u, f \| \alpha)=0 \text { for every } \alpha \in G \text { and every } h\right.  \tag{15.4}\\
\text { such that } \operatorname{det}(h)=0\} \\
\mathcal{T}_{\omega}^{p}=\left\{f \in \mathcal{R}_{\omega}^{p} \mid\left\langle f, \mathcal{S}_{\omega}\right\rangle=0\right\}  \tag{15.5}\\
\mathcal{R}_{\omega}^{p}(\Gamma)=\mathcal{R}_{\omega}^{p} \cap \mathcal{N}_{\omega}^{p}(\Gamma), \quad \mathcal{T}_{\omega}^{p}(\Gamma)=\mathcal{T}_{\omega}^{p} \cap \mathcal{N}_{\omega}^{p}(\Gamma) \tag{15.6}
\end{gather*}
$$

Here $\omega$ is either a representation or a quasi-representation of $\mathfrak{K}$, and $\Gamma$ is a congruence subgroup of $G=U\left(\eta_{n}\right)$ or $\mathcal{G}$ (see $\S 14.13$ ). Clearly $\mathcal{S}_{\omega}=\mathcal{M}_{\omega} \cap \mathcal{R}_{\omega}^{p}$. Notice that $\langle f, g\rangle$ is meaningful for $f, g \in \mathcal{R}_{\omega}^{p}$, since the elements of $\mathcal{R}_{\omega}^{p}$ are rapidly decreasing at the cusps.
15.6. Proposition (Cases SP and UT). The notation being as above, the following assertions hold:
(1) $\mathcal{T}_{\omega}^{p}(\Gamma)=\left\{f \in \mathcal{R}_{\omega}^{p}(\Gamma) \mid\left\langle f, \mathcal{S}_{\omega}(\Gamma)\right\rangle=0\right\}$.
(2) $\mathcal{R}_{\omega}^{p}(\Gamma)=\mathcal{S}_{\omega}(\Gamma) \oplus \mathcal{T}_{\omega}^{p}(\Gamma)$ and $\mathcal{R}_{\omega}^{p}=\mathcal{S}_{\omega} \oplus \mathcal{T}_{\omega}^{p}$.
(3) Let $\mathfrak{p}: \mathcal{R}_{\omega}^{p} \rightarrow \mathcal{S}_{\omega}$ be the projection map obtained from the last direct sum decomposition of (2). Then $\left(\mathcal{R}_{\omega}^{p}\right)^{\varepsilon}=\mathcal{R}_{\omega^{\varepsilon}}^{p \varepsilon}$ and $\mathfrak{p}(f)^{\varepsilon}=\mathfrak{p}\left(f^{\varepsilon}\right)$ for every $\varepsilon \in \operatorname{Aut}(\mathbf{C})$, and $\langle f, g\rangle=\langle\mathfrak{p}(f), g\rangle$ for every $f \in \mathcal{R}_{\omega}^{p}$ and every $g \in \mathcal{S}_{\omega}$.

Proof. Clearly the left-hand side of the equality of (1) is contained in the right-hand side. Let $f$ be an element of the right-hand side and let $g \in \mathcal{S}_{\omega}$. Take a subgroup $\Gamma^{\prime}$ of $\Gamma$ of finite index so that $g \in \mathcal{S}_{\omega}\left(\Gamma^{\prime}\right)$ and take a set of representatives $A$ for $\Gamma^{\prime} \backslash \Gamma$. Then $\sum_{\alpha \in A} g \| \alpha \in \mathcal{S}_{\omega}(\Gamma)$, so that $0=\sum_{\alpha \in A}\langle f, g \| \alpha\rangle=$ $\sum_{\alpha \in A}\left\langle f \| \alpha^{-1}, g\right\rangle=\left[\Gamma: \Gamma^{\prime}\right]\langle f, g\rangle$. Thus $f \in \mathcal{T}_{\omega}^{p}(\Gamma)$. This proves (1). By Lemma $14.3, \mathcal{R}_{\omega}^{p}(\Gamma)$ is finite-dimensional over $\mathbf{C}$, and hence from (1) we obtain the first equality of (2), which clearly implies the second equality. To prove (3), take $f \in \mathcal{N}_{\omega}^{p}$; modifying the expression for $q f$ in the proof of Theorem 14.9, we can put $q f=\sum_{s}(\pi i)^{-|s|} \theta_{X}^{s} D_{\zeta \otimes \sigma^{s}}^{s} g_{s}$ with $g_{s} \in \mathcal{M}_{\zeta \otimes \sigma^{s}}$. By (14.16) we have $q^{\varepsilon} f^{\varepsilon}=(q f)^{\varepsilon}=\sum_{s}(\pi i)^{-|s|} \theta_{X}^{s \varepsilon} D_{\zeta^{\varepsilon} \otimes \sigma^{s \varepsilon}}^{s \varepsilon} \ell_{s}$. If $\omega$ is a representation of $\mathfrak{K}$, then given $\alpha \in G$, take $\beta$ as in Lemma 10.5. Then by that lemma and Proposition 12.10 (2), $(q f)^{\varepsilon} \|_{\varsigma^{\varepsilon}} \alpha=\left((q f) \|_{\zeta} \beta\right)^{\varepsilon}$ and $q^{\varepsilon} \|_{\kappa \mathbf{b}} \alpha=\left(q \|_{\kappa \mathbf{b}} \beta\right)^{\varepsilon}$, and hence $f^{\varepsilon} \|_{\omega^{\varepsilon}} \alpha=\left(f \|_{\omega} \beta\right)^{\varepsilon}$. Suppose $f \in \mathcal{R}_{\omega}^{p}$; then from the last equality we see that $f^{\varepsilon} \in \mathcal{R}_{\omega^{\varepsilon}}^{p}$. This shows that $\left(\mathcal{R}_{\omega}^{p}\right)^{\varepsilon} \subset \mathcal{R}_{\omega^{\varepsilon}}^{p}$. Considering the action of $\varepsilon^{-1}$ similarly, we obtain $\left(\mathcal{R}_{\omega}^{p}\right)^{\varepsilon}=\mathcal{R}_{\omega^{\varepsilon}}^{p}$. If $\omega$ is a quasi-representation, given $\alpha \in \mathcal{G}$, we can find $\beta \in \mathcal{G}$ with which we can make a similar type of argument, as shown in the proof of Theorem 10.8. Then we obtain the desired result. Let $\Lambda=\sum_{v \in \mathbf{a}} L_{\omega, v}$ and $\Lambda^{\prime}=\sum_{v \in \mathbf{a}} L_{\omega^{\varepsilon}, v}$ (see Corollary 12.16). From the proof of Theorem 14.12 we easily see that $\Lambda$ maps $\mathcal{R}_{\omega}^{p}$ into itself. Then Theorem 12.15 shows that $\Lambda$ is a hermitian operator on $\mathcal{R}_{\omega}^{p}(\Gamma)$, so that $\mathcal{R}_{\omega}^{p}(\Gamma)$ is a finite direct sum $\bigoplus_{\mu} \mathcal{E}_{\mu}$, where $\mathcal{E}_{\mu}=\left\{f \in \mathcal{R}_{\omega}^{p}(\Gamma) \mid \Lambda f=\mu f\right\}$. By Corollary 12.16 we have $\mathcal{S}_{\omega}(\Gamma)=\mathcal{E}_{0}$, so that $\mathcal{T}_{\omega}^{p}(\Gamma)=\sum_{\mu \neq 0} \mathcal{E}_{\mu}$ by (1). Let $f \in \mathcal{R}_{\omega}^{p}(\Gamma)$; then $f=g+\sum_{\mu \neq 0} h_{\mu}$ with $g \in \mathcal{S}_{\omega}(\Gamma)$ and $h_{\mu} \in \mathcal{E}_{\mu}$, and $f^{\varepsilon}=g^{\varepsilon}+\sum_{\mu \neq 0} h_{\mu}^{\varepsilon}$. By Theorem 10.8 (1), $g^{\varepsilon} \in \mathcal{S}_{\omega^{\varepsilon}}$. Now, for a fixed $h \in \mathcal{N}_{\omega}^{p}$, take $\pi i E_{v} h$ and $\omega$ as $g_{s}$ and $\zeta$ in (14.16). By Theorem 14.12 (3) we have $\ell_{s}=\pi i E_{v \varepsilon} h^{\varepsilon}$, and hence (14.16) shows that $\left(L_{\omega, v} h\right)^{\varepsilon}=L_{\omega^{\varepsilon}, v \varepsilon} h^{\varepsilon}$. Thus $(\Lambda h)^{\varepsilon}=\Lambda^{\prime} h^{\varepsilon}$, so that we have $\Lambda^{\prime} h_{\mu}^{\varepsilon}=\mu^{\varepsilon} h_{\mu}^{\varepsilon}$ and hence $\left\langle\mathcal{S}_{\omega^{\varepsilon}}, h_{\mu}^{\varepsilon}\right\rangle=0$. Thus $\sum_{\mu \neq 0} h_{\mu}^{\varepsilon} \in \mathcal{T}_{\omega^{\varepsilon}}^{p}$. This shows that $g^{\varepsilon}=\mathfrak{p}\left(f^{\varepsilon}\right)$, which is the second equality of (3). Since $f-\mathfrak{p}(f) \in \mathcal{T}_{\omega}^{p}$, the last equality of (3) follows from (15.5).
15.7. Proposition (Case UB). Suppose that $\Gamma \backslash \mathcal{H}$ is compact for a congruence subgroup $\Gamma$ of $G$. Put $\mathcal{T}_{\omega}^{p}=\left\{f \in \mathcal{N}_{\omega}^{p} \mid\left\langle f, \mathcal{S}_{\omega}\right\rangle=0\right\}$ and $\mathcal{T}_{\omega}^{p}(\overline{\mathbf{Q}})=\mathcal{T}_{\omega}^{p} \cap \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})$. Then the following assertions hold:
(1) $\mathcal{T}_{\omega}^{p}(\Gamma)=\left\{f \in \mathcal{N}_{\omega}^{p}(\Gamma) \mid\left\langle f, \mathcal{S}_{\omega}(\Gamma)\right\rangle=0\right\}$.
(2) $\mathcal{N}_{\omega}^{p}(\Gamma)=\mathcal{S}_{\omega}(\Gamma) \oplus \mathcal{T}_{\omega}^{p}(\Gamma)$ and $\mathcal{N}_{\omega}^{p}=\mathcal{S}_{\omega} \oplus \mathcal{T}_{\omega}^{p}$.
(3) $\mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})=\mathcal{S}_{\omega}(\overline{\mathbf{Q}}) \oplus \mathcal{T}_{\omega}^{p}(\overline{\mathbf{Q}})$ and $\mathcal{T}_{\omega}^{p}=\mathcal{T}_{\omega}^{p}(\overline{\mathbf{Q}}) \otimes_{\overline{\mathbf{Q}}} \mathbf{C}$.

Proof. Under our assumption, $\langle f, g\rangle$ is meaningful for every $f, g \in \mathcal{N}_{\omega}^{p}$, and
so (1) and (2) can be proved in the same manner as for Proposition 15.6. To prove (3), put $\Lambda=\sum_{v \in \mathbf{a}^{\prime}} L_{\omega, v}$. Then, for the same reason as in the proof of Proposition 15.6, $\mathcal{N}_{\omega}^{p}(\Gamma)=\bigoplus_{\mu} \mathcal{E}_{\mu}$ with $\mathcal{E}_{\mu}=\left\{f \in \mathcal{N}_{\omega}^{p}(\Gamma) \mid \Lambda f=\mu f\right\}, \mathcal{S}_{\omega}(\Gamma)=\mathcal{E}_{0}$, and $\mathcal{T}_{\omega}^{p}(\Gamma)=$ $\sum_{\mu \neq 0} \mathcal{E}_{\mu}$. By Theorem 14.9 (2) and Lemma 14.6 (iv), $\Lambda$ maps $\mathcal{N}_{\omega}^{p}(\Gamma, \overline{\mathbf{Q}})$ into itself. Therefore, in view of the last equality of Proposition 14.10 we see that $\mu \in \overline{\mathbf{Q}}$ and $\mathcal{E}_{\mu}=\left(\mathcal{E}_{\mu} \cap \mathcal{N}_{\omega}^{p}(\overline{\mathbf{Q}})\right) \otimes_{\overline{\mathbf{Q}}} \mathbf{C}$, from which we obtain (3) in a straightforward way.
15.8. Lemma (Cases SP and UT). Let $W$ be a subfield of $\mathbf{C}$ containing the Galois closure of $K$ in $\mathbf{C}$ over $\mathbf{Q}$; let $g \in \mathcal{M}_{l}(W)$ and $h \in \mathcal{M}_{l^{\prime}}(W)$ with $l, l^{\prime} \in$ $2^{-1} \mathbf{Z}^{\mathbf{a}} ;$ further define the operator $\Delta_{l^{\prime}}^{p}$ by $\Delta_{l^{\prime}}^{p} h=\left(D_{\rho^{\prime}}^{Z} h\right)(\psi)$ with $\rho^{\prime}(a, b)=\operatorname{det}(b)^{l^{\prime}}$, $Z=\bigotimes_{v \in \mathbf{a}} Z_{v}=\mathbf{C} \psi \subset S_{n p}(T)$, where $Z_{v}=\mathbf{C} \psi_{v} \subset S_{n p_{v}}\left(T_{v}\right), \psi=\prod_{v \in \mathbf{a}} \psi_{v}$, and $\psi_{v}=\operatorname{det}(x)^{p_{v}}$. Put $l_{0}=\operatorname{Min}\left\{l_{v}, l_{v}^{\prime} \mid v \in \mathbf{a}\right\}$ and $k=l+l^{\prime}+2 p$. Suppose that $l_{0} \geq n / 2$ in Case $S P$ and $l_{0} \geq n$ in Case UT. Then there exists an element $q$ of $\mathcal{M}_{k}(W)$ such that $(\pi i)^{-n|p|}\left\langle g \Delta_{l^{\prime}}^{p} h, f\right\rangle=\langle q, f\rangle$ for every $f \in \mathcal{S}_{k}$.

Proof. From (12.24c) and Theorem 14.12 (4) we see that $(\pi i)^{-n|p|} \Delta_{l^{\prime}}^{p}$ sends $\mathcal{M}_{l^{\prime}}(W)$ into $\mathcal{N}_{l^{\prime}+2 p}^{n p}(W)$. Therefore, if $g$ or $h$ is a cusp form, then $g \Delta_{l^{\prime}}^{p} h$ belongs to $\mathcal{R}_{\omega}^{n p}$ with $\omega(a, b)=\operatorname{det}(b)^{k}$, and hence our assertion follows from Proposition 15.6 (3). If neither $g$ nor $h$ is a cusp form, then $l=\mu \mathbf{a}$ and $l^{\prime}=\nu \mathbf{a}$ with $\mu, \nu \in 2^{-1} \mathbf{Z}$ by virtue of [S97, Proposition 10.6 (3)]. This case will be proven in §A8.8.

If $G=S L_{2}(F)$, we can prove a result stronger than Proposition 15.3 and Proposition 15.6. In fact, for every weight $k \in 2^{-1} \mathbf{Z}^{\mathbf{a}}$, there exists a C-linear map $\mathfrak{p}_{k}: \bigcup_{p} \mathcal{N}_{k}^{p} \rightarrow \mathcal{S}_{k}$ with the following properties:
(15.7a) $\langle f, h\rangle=\left\langle\mathfrak{p}_{k}(f), h\right\rangle$ for every $f \in \bigcup_{p} \mathcal{N}_{k}^{p}$ and $h \in \mathcal{S}_{k}$;
(15.7b) $\mathfrak{p}_{k}(f)^{\sigma}=\mathfrak{p}_{k \sigma}\left(f^{\sigma}\right)$ for every $\sigma \in \operatorname{Aut}(\mathbf{C})$.

See [S87b, Proposition 9.4] for the proof.

## CHAPTER IV

## EISENSTEIN SERIES OF SIMPLER TYPES

## 16. Eisenstein series on $U\left(\eta_{n}\right)$

16.1. This section concerns Cases SP and UT. Thus $G=S p(n, F)$ or $G=U\left(\eta_{n}\right)$ as in $\S 3.5$, and $G_{1}=G \cap S L_{2 n}(K)$. We retain the convention that $K=F, \mathfrak{r}=\mathfrak{g}$, and $\rho=\operatorname{id}_{F}$ in Case SP. We start with preliminary discussions on certain infinite series that appear as nonarchimedean factors of a Fourier coefficient of an Eisenstein series. We first put

$$
\begin{gather*}
S=\left\{h \in K_{n}^{n} \mid h^{*}=h\right\},  \tag{16.1a}\\
S(\mathfrak{a})=S \cap(\mathfrak{r a})_{n}^{n},  \tag{16.1b}\\
S_{\mathbf{h}}(\mathfrak{a})=\prod_{v \in \mathbf{h}} S(\mathfrak{a})_{v}, \\
\widetilde{S}=\widetilde{S}^{n}=S_{\mathbf{a}} \times \prod_{v \in \mathbf{h}} \widetilde{S}_{v}\left(\subset S_{\mathbf{A}}\right), \quad \widetilde{S}_{v}=\widetilde{S}_{v}^{n}=\left\{x \in S_{v} \mid \operatorname{tr}\left(x \cdot S(\mathfrak{r})_{v}\right) \subset \mathfrak{g}_{v}\right\} .
\end{gather*}
$$

Here $\mathfrak{a}$ is a fractional ideal in $F$ or $K$. We need the symbols $\mathbf{e}, \mathbf{e}_{\mathbf{h}}, \mathbf{e}_{\mathbf{a}}$, and $\mathbf{e}_{\mathbf{A}}$ introduced in §1.6; we also recall that $\mathbf{e}_{\mathbf{a}}^{n}(X)=\exp \left(2 \pi i \sum_{v \in \mathbf{a}} \operatorname{tr}(X)\right)$ for $X \in\left(\mathbf{C}_{n}^{n}\right)^{\mathbf{a}}$ and $\mathbf{e}_{\mathbf{a}}(x)=\mathbf{e}_{\mathbf{a}}^{1}(x)$ for $x \in \mathbf{C}^{\mathbf{a}}$, as defined in (5.15) and (5.16). Similarly we put

$$
\begin{equation*}
\mathbf{e}_{\mathbf{x}}^{n}(W)=\mathbf{e}_{\mathbf{x}}(\operatorname{tr}(W)) \text { for } W \in\left(K_{\mathbf{x}}\right)_{n}^{n} \text { such that } \operatorname{tr}(W) \in F_{\mathbf{x}}, \tag{16.2}
\end{equation*}
$$

where $\mathbf{x}$ is one of the symbols $v, \mathbf{h}$, and $\mathbf{A}$. For example, $\mathbf{e}_{v}^{n}(Y Z)$ is meaningful for $Y, Z \in S_{v}$ in both Cases SP and UT.

For $\sigma \in S_{\mathbf{A}}$ we define an integral $\mathfrak{r}$-ideal $\nu_{0}(\sigma)$ and a positive integer $\nu(\sigma)$ by (1.20) and (1.21). Since $\sigma^{*}=\sigma$, we can show (see [S97, §13.4]) that

$$
\begin{equation*}
\nu_{0}(\sigma)=\left(\mathfrak{g} \cap \nu_{0}(\sigma)\right) \mathfrak{r} . \tag{16.3}
\end{equation*}
$$

We then put

$$
\begin{equation*}
\nu[\sigma]=N\left(\mathfrak{g} \cap \nu_{0}(\sigma)\right) \quad\left(\sigma \in S_{\mathbf{A}}\right) . \tag{16.4}
\end{equation*}
$$

Notice that $\nu(\sigma)=\nu[\sigma]^{[K: F]}$. For $s \in S_{\mathrm{A}}$ in Case SP we put

$$
\begin{gather*}
\gamma(s)=\prod_{v \in \mathbf{h}} \gamma_{v}(s), \quad \gamma_{v}(s)=\int_{L_{v}} \mathbf{e}_{v}\left(x s \cdot{ }^{t} x / 2\right) d x, \quad L_{v}=\left(\mathfrak{g}_{v}\right)_{n}^{1}  \tag{16.5}\\
\omega(s)=\gamma(s) /|\gamma(s)| \tag{16.6}
\end{gather*}
$$

where $d x$ is the Haar measure of $\left(F_{v}\right)_{n}^{1}$ such that $\int_{L_{v}} d x=1$, and we assume that $\gamma(s) \neq 0$ in (16.6). Clearly $\gamma_{v}(s)=1$ for almost all $v$, and so the product over all $v \in \mathbf{h}$ is meaningful. Since we view $S_{v}$ as a subset of $S_{\mathbf{A}}$, we can speak of $\gamma(s)$ and $\omega(s)$ for $s \in S_{v}$. We shall show in Lemma A1.6 that $\gamma_{v}(s) \neq 0$ if $v \nmid 2$.

Let $\mathfrak{d}$ be the different of $F$ relative to $\mathbf{Q}$. Take an element $\delta$ of $F_{\mathrm{h}}^{\times}$such that $\mathfrak{d}=\delta \mathfrak{g}$. Given $\zeta \in \widetilde{S}$, we put

$$
\begin{gather*}
\alpha_{\mathfrak{c}}^{0}(\zeta, s)=\prod_{v \nmid c} \alpha_{v}(\zeta, s), \quad \alpha_{v}^{0}(\zeta, s)=\sum_{\sigma \in S_{v} / S(\mathfrak{r})_{v}} \mathbf{e}_{v}^{n}\left(-\delta_{v}^{-1} \zeta \sigma\right) \nu[\sigma]^{-s},  \tag{16.7a}\\
\alpha_{\mathfrak{c}}^{1}(\zeta, s)=\prod_{v \nmid c} \alpha_{v}^{1}(\zeta, s), \quad \alpha_{v}^{1}(\zeta, s)=\sum_{\sigma \in S_{v} / S(\mathfrak{g})_{v}} \omega\left(\delta_{v}^{-1} \sigma\right) \mathbf{e}_{v}^{n}\left(-\delta_{v}^{-1} \zeta \sigma\right) \nu(\sigma)^{-s} . \tag{16.7b}
\end{gather*}
$$

The series $\alpha_{c}^{0}$ and $\alpha_{v}^{0}$ are defined in both Cases SP and UT, but $\alpha_{c}^{1}$ and $\alpha_{v}^{1}$ are defined only in Case SP under the condition that $v \nmid 2$ and $\mathfrak{c} \subset 4 \mathfrak{g}$. Clearly $\gamma_{v}(s)=\gamma_{v}(s+b)$ if $b \prec\left(2 / \delta_{v}\right) \mathfrak{g}_{v}$, and hence $\omega_{v}(s)=\omega_{v}(s+b)$ for such a $b$. Since $\nu_{0}(\sigma)$ depends only on $\sigma$ modulo $\left(\mathfrak{r}_{v}\right)_{n}^{n}$, we see that the series of (16.7a, b) are formally well-defined. These series appear as nonarchimedean factors of the Fourier coefficients of Eisenstein series, as will be shown in Proposition 16.9 below.

We call an element $\psi$ of $\widetilde{S}_{v}$ regular if:

$$
\begin{aligned}
& K=F, n \text { is even, and } \operatorname{det}(2 \psi) \in \mathfrak{g}_{v}^{\times} ; \text {or } \\
& K=F, n \text { is odd, and } \operatorname{det}(2 \psi) \in 2 \mathfrak{g}_{v}^{\times} ; \text {or } \\
& K \neq F, n \text { is even, and } \operatorname{det}(\varepsilon \psi) \in \mathfrak{r}_{v}^{\times}, \text {where } \varepsilon \text { is an element of } K_{v}^{\times} \\
& \quad \text { such that } \varepsilon^{-1} \mathfrak{r}_{v}=\left\{x \in K_{v} \mid \operatorname{Tr}_{K / F}\left(x \mathfrak{r}_{v}\right) \subset \mathfrak{g}_{v}\right\} ; \text { or }
\end{aligned}
$$

$K \neq F, v$ is unramified in $K, n$ is odd, and $\operatorname{det}(\psi) \in \mathfrak{g}_{v}^{\times}$.
In Case SP, for $\xi \in \widetilde{S}_{v}^{r} \cap G L_{r}\left(F_{v}\right)$ we define $\lambda(\xi)$ as follows: Put $h=(-1)^{r / 2} \operatorname{det}(\xi)$ if $r$ is even, and $h=2(-1)^{(r-1) / 2} \operatorname{det}(\xi)$ if $r$ is odd; then $\lambda(\xi)=1$ if $F_{v}\left(h^{1 / 2}\right)=$ $F_{v}, \lambda(\xi)=-1$ if $F_{v}\left(h^{1 / 2}\right)$ is an unramified quadratic extension of $F_{v}$, and $\lambda(\xi)=0$ if $F_{v}\left(h^{1 / 2}\right)$ is ramified over $F_{v}$.
16.2. Theorem. Let $\zeta \in \widetilde{S}_{v}^{n}, v \in \mathbf{h}$, and $\operatorname{rank}(\zeta)=r$; suppose $\zeta=0$ or $\zeta=\operatorname{diag}[\xi, 0]$ with $\xi \in \widetilde{S}_{v}^{r} \cap G L_{r}\left(K_{v}\right)$; put $q=\left|\pi_{v}\right|^{-1}$ with a prime element $\pi_{v}$ of $F_{v}$. Define power series $A_{\zeta}^{0}(t)$ and $A_{\zeta}^{1}(t)$ in an indeterminate $t$ by

$$
\begin{aligned}
A_{\zeta}^{0}(t) & =\sum_{\sigma \in S_{v} / S(\mathfrak{r})_{v}} \mathbf{e}_{v}^{n}\left(-\delta_{v}^{-1} \zeta \sigma\right) t^{e(\sigma)} \\
A_{\zeta}^{1}(t) & =\sum_{\sigma \in S_{v} / S(\mathfrak{g})_{v}} \omega\left(\delta_{v}^{-1} \sigma\right) \mathbf{e}_{v}^{n}\left(-\delta_{v}^{-1} \zeta \sigma\right) t^{e(\sigma)}
\end{aligned}
$$

where $e(\sigma)$ is the integer defined by $\nu[\sigma]=q^{e(\sigma)}$. (This means that $A_{\zeta}^{i}\left(q^{-s}\right)=$ $\alpha_{v}^{i}(\zeta, s)$.) Then $A_{\zeta}^{0}=f_{\zeta}^{0} g_{\zeta}^{0}$ and $A_{\zeta}^{1}(t)=f_{\zeta}^{1}(t) g_{\zeta}^{1}\left(q^{-1 / 2} t\right)$ with polynomials $g_{\zeta}^{i}$ with coefficients in $\mathbf{Z}$ whose constant terms are 1 and rational functions $f_{\zeta}^{i}$ given as follows:

$$
\begin{aligned}
f_{\zeta}^{0}(t) & =\frac{(1-t) \prod_{i=1}^{[n / 2]}\left(1-q^{2 i} t^{2}\right)}{\left(1-\lambda q^{(2 n-r) / 2} t\right) \prod_{i=1}^{[(n-r) / 2]}\left(1-q^{2 n-r-2 i+1} t^{2}\right)} \quad(\text { Case SP, } r \in 2 \mathbf{Z}), \\
f_{\zeta}^{0}(t) & =\frac{(1-t) \prod_{i=1}^{[n / 2]}\left(1-q^{2 i} t^{2}\right)}{\prod_{i=1}^{[(n-r+1) / 2]}\left(1-q^{2 n-r-2 i+2} t^{2}\right)} \quad(\text { Case SP, } r \notin 2 \mathbf{Z}), \\
f_{\zeta}^{0}(t) & =\frac{\prod_{i=1}^{n}\left(1-\tau^{i-1} q^{i-1} t\right)}{\prod_{i=1}^{n-r}\left(1-\tau^{n+i} q^{n+i-1} t\right)} \quad \quad(\text { Case UT }), \\
f_{\zeta}^{1}(t) & =\frac{\prod_{i=1}^{[(n+1) / 2]}\left(1-q^{2 i-1} t^{2}\right)}{\left(1-\lambda q^{(2 n-r) / 2} t\right) \prod_{i=1}^{[(n-r) / 2]}\left(1-q^{2 n+1-r-2 i} t^{2}\right)} \quad(r \notin 2 \mathbf{Z}),
\end{aligned}
$$

$$
f_{\zeta}^{1}(t)=\frac{\prod_{i=1}^{[(n+1) / 2]}\left(1-q^{2 i-1} t^{2}\right)}{\prod_{i=1}^{[(n-r+1) / 2]}\left(1-q^{2 n+2-r-2 i} t^{2}\right)} \quad(r \in 2 \mathbf{Z})
$$

Here $\lambda=1$ if $\zeta=0$ and $\lambda=\lambda(\xi)$ otherwise; $\tau^{i}$ is the symbol defined as follows:

$$
\tau^{i}=\left\{\begin{array}{l}
1 \quad \text { if } i \text { is even or } v \text { splits in } K,  \tag{16.8}\\
-1 \quad \text { if } i \text { is odd and } v \text { remains prime in } K, \\
0 \quad \text { if } i \text { is odd and } v \text { is ramified in } K .
\end{array}\right.
$$

Moreover, $A_{\zeta}^{0}=f_{\zeta}^{0}$ if $\xi$ is regular or $\zeta=0$, except when $K \neq F, v$ is ramified in $K$, and $\zeta \neq 0 ; A_{\zeta}^{1}=f_{\zeta}^{1}$ if $\xi \in G L_{r}\left(\mathfrak{g}_{v}\right)$ or $\zeta=0$.

The assertion concerning $A_{\zeta}^{0}$ and the formulas for $f_{\zeta}^{0}$ were given in [S97, Theorem 13.6]. We shall prove the part concerning $A_{\zeta}^{1}$ and $f_{\zeta}^{1}$ in §A1.9.
16.3. By a Hecke character of an algebraic number field $K$ we understand a continuous homomorphism $\psi$ of $K_{\mathbf{A}}^{\times}$into $\mathbf{T}$ such that $\psi\left(K^{\times}\right)=1$. For such a $\psi$ we denote by $\psi_{v}, \psi_{\mathbf{a}}$, and $\psi_{\mathbf{h}}$ its restrictions to $K_{v}^{\times}, K_{\mathbf{a}}^{\times}$, and $K_{\mathbf{h}}^{\times}$, respectively. Also we denote by $\psi^{*}$ the ideal character associated with $\psi$. We put $\psi^{*}(\mathfrak{a})=0$ if $\mathfrak{a}$ is not prime to the conductor of $\psi$. In the setting of $\S 16.1$, we have $\psi_{\mathbf{a}}(x)=x_{\mathbf{a}}^{h}\left|x_{\mathbf{a}}\right|^{i \kappa-h}$ with $h \in \mathbf{Z}^{\mathbf{a}}$ and $\kappa \in \mathbf{R}^{\mathbf{a}}$. In this book we always assume that $\psi$ is normalized in the sense that $\sum_{v \in \mathbf{a}} \kappa_{v}=0$.

Given a Hecke character $\psi$ of $K$, we define the $L$-function $L(s, \psi)$ as usual, and its partial series $L_{\mathfrak{c}}(s, \psi)$ for an integral ideal $\mathfrak{c}$ in $F$ by

$$
\begin{equation*}
L_{\mathfrak{c}}(s, \psi)=\prod_{\mathfrak{p} \nmid \mathfrak{c}}\left[1-\psi^{*}(\mathfrak{p}) N(\mathfrak{p})^{-s}\right]^{-1}=L(s, \psi) \prod_{\mathfrak{p} \mid \mathfrak{c}}\left[1-\psi^{*}(\mathfrak{p}) N(\mathfrak{p})^{-s}\right] . \tag{16.9}
\end{equation*}
$$

Here $\mathfrak{p}$ denotes a prime ideal in $K$. We also put $\psi_{\mathfrak{c}}=\prod_{v \mid \mathfrak{c}} \psi_{v}$.
16.4. For $x=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right] \in G L_{2 n}(K)_{\mathbf{A}}$ with $a \in\left(K_{\mathbf{A}}\right)_{n}^{n}$, we write $a=a_{x}, b=$ $b_{x}, c=c_{x}$, and $d=d_{x}$. We let $G_{\mathbf{A}}$ act on $\mathcal{H}$ through the projection map $G_{\mathbf{A}} \rightarrow G_{\mathbf{a}}$ (see $\S 3.5$ ), and define $j_{\alpha}(z)=j(\alpha, z)$ for $\alpha \in G_{\mathbf{A}}$ and $z \in \mathcal{H}$ by (5.3). We put

$$
\begin{gather*}
P=\left\{\xi \in G \mid c_{\xi}=0\right\}  \tag{16.10}\\
Q=\left\{\xi \in P \mid b_{\xi}=0\right\}, \quad R=\left\{\xi \in P \mid a_{\xi}=1\right\} \tag{16.11}
\end{gather*}
$$

We already introduced modular forms of half-integral weight in §6.10. To deal with them adelically we need to consider the metaplectic group $M_{\mathbf{A}}$ and its subset $\mathfrak{M}$ defined by

$$
\begin{equation*}
\mathfrak{M}=\left\{\sigma \in M_{\mathbf{A}} \mid \operatorname{pr}(\sigma) \in P_{\mathbf{A}} C^{\theta}\right\} \tag{16.12}
\end{equation*}
$$

where pr is the map of (16.13) below and $C^{\theta}$ is defined by (6.29) (see also (A2.12a, b)). We need the following facts $(1 \sim 4)$ which will be explained in detail in $\S A 2$ in the Appendix I:
(16.13) There is a surjective homomorphism $\mathrm{pr}: M_{\mathbf{A}} \rightarrow G_{\mathbf{A}}=\operatorname{Sp}\left(n, F_{\mathbf{A}}\right)$ whose kernel is $\mathbf{T}$ (viewed as a subgroup of $M_{\mathbf{A}}$ in a certain way).
(16.14) There is a lift $G \rightarrow M_{\mathbf{A}}$ which combined with pr gives the identity map on $G$.
(16.15) There is also a lift $r_{P}: P_{\mathbf{A}} \rightarrow M_{\mathbf{A}}$ in the same sense that coincides with the lift of (16.14) on $P$.

For $\alpha=\operatorname{pr}(\sigma)$ with $\sigma \in M_{\mathbf{A}}$ and $z \in \mathcal{H}$ we put $a_{\sigma}=a_{\alpha}, b_{\sigma}=b_{\alpha}, c_{\sigma}=c_{\alpha}, d_{\sigma}=d_{\alpha}$, $\sigma(z)=\sigma z=\alpha z$, and $j_{\sigma}(z)=j_{\alpha}(z)$.
(16.16) For every $\sigma \in \mathfrak{M}$ there is a holomorphic function $h_{\sigma}(z)=h(\sigma, z)$ of $z \in \mathcal{H}$ with the following properties:
(16.16a) $h(\sigma, z)^{2}=\zeta \cdot j(\operatorname{pr}(\sigma), z)^{\mathbf{a}}$ with $\zeta \in \mathbf{T} ; h(\sigma, z) \in \mathbf{T}$ if $\operatorname{pr}(\sigma)_{\mathbf{a}}=1$;
(16.16b) $h\left(t \cdot r_{P}(\gamma), z\right)=t^{-1}\left|\operatorname{det}\left(d_{\gamma}\right)_{\mathbf{a}}\right|_{\mathbf{A}}^{1 / 2} \quad$ if $t \in \mathbf{T} \quad$ and $\quad \gamma \in P_{\mathbf{A}}$;
(16.16c) $h(\rho \sigma \tau, z)=h(\rho, z) h(\sigma, \tau z) h(\tau, z)$ if $\operatorname{pr}(\rho) \in P_{\mathbf{A}}$ and $\operatorname{pr}(\tau) \in C^{\theta}$;
(16.16d) $h(\sigma, z)$ coincides with $h$ of Theorem 6.8 if $\sigma$ belongs to $\Gamma^{\theta}$ of (6.30).

We shall always view $G$ and its subgroups as subgroups of $M_{\mathbf{A}}$ by means of the lift of (16.14). In $\S 14.14$ we defined a group $\mathcal{G}$ consisting of all $(\alpha, p)$ such that $p^{2} / j_{\alpha}^{\text {a }}$ is a root of unity. We can view $\mathcal{G}$ as a subgroup of $M_{\mathbf{A}}$ by identifying $(\alpha, p)$ with the element $\sigma$ of $M_{\mathbf{A}}$ such that $\operatorname{pr}(\sigma)=\alpha_{\mathbf{a}}$ and $h_{\sigma}=p$.

Let $k$ be an integral or a half-integral weight. Here an integral weight means an element of $\mathbf{Z}^{\mathbf{b}}$; a half-integral weight occurs only in Case SP; see $\S 6.10$. We define $\mathcal{M}_{k}$ as in $\S \S 5.5$ and 6.10 . Now we define a factor of automorphy $j^{k}$ by

$$
j_{\sigma}^{k}(z)=j^{k}(\sigma, z)=\left\{\begin{array}{lr}
j(\sigma, z)^{k} & \left(\sigma \in G_{\mathbf{A}}, k \in \mathbf{Z}^{\mathbf{b}}\right)  \tag{16.17}\\
h_{\sigma}(z) j_{\sigma}(z)^{[k]} & \left(\sigma \in \mathfrak{M}, k \notin \mathbf{Z}^{\mathbf{b}}\right)
\end{array}\right.
$$

where $[k]=\left(k_{v}-1 / 2\right)_{v \in \mathbf{a}}$, and we are employing the notation of (5.4a). We put $[k]=k$ if $k \in \mathbf{Z}^{\mathbf{b}}$. Then, given $f: \mathcal{H} \rightarrow \mathbf{C}$ and $\sigma \in G_{\mathbf{A}}$ or $\sigma \in \mathfrak{M}$, we define $f \|_{k} \sigma: \mathcal{H} \rightarrow \mathbf{C}$ by

$$
\begin{equation*}
\left(f \|_{k} \sigma\right)(z)=j_{\sigma}^{k}(z)^{-1} f(\sigma z) \tag{16.18}
\end{equation*}
$$

For $k \in \mathbf{Z}^{\mathbf{b}}$ this is consistent with (5.7). If $k$ is half-integral, $\sigma \in G \cap \mathfrak{M}$, and we identify $\sigma$ with the element $\left(\sigma, h_{\sigma}\right)$ of $\mathcal{G}$ of $\S 14.14$, then (16.18) is consistent with (14.18b); notice that $f\left\|_{k}(\sigma \tau)=\left(f \|_{k} \sigma\right)\right\|_{k} \tau$ if $\operatorname{pr}(\sigma) \in P_{\mathbf{A}}$ or $\operatorname{pr}(\tau) \in C^{\theta}$. From (16.16b) and (16.17) we obtain

$$
\begin{equation*}
j^{k}\left(r_{P}(p), z\right)=\left|\operatorname{det}\left(d_{p}\right)\right|^{k-[k]} \operatorname{det}\left(d_{p}\right)^{[k]} \quad \text { if } \quad p \in P_{\mathbf{A}} . \tag{16.19}
\end{equation*}
$$

16.5. For two fractional ideals $\mathfrak{x}$ and $\mathfrak{y}$ in $F$ such that $\mathfrak{x y} \subset \mathfrak{g}$ we put

$$
\begin{gather*}
D[\mathfrak{x}, \mathfrak{y}]=\left\{x \in G_{\mathbf{A}} \mid a_{x} \prec \mathfrak{r}, b_{x} \prec \mathfrak{r x}, c_{x} \prec \mathfrak{r y}, d_{x} \prec \mathfrak{r}\right\},  \tag{16.20a}\\
D_{0}[\mathfrak{r}, \mathfrak{y}]=\left\{x \in D[\mathfrak{r}, \mathfrak{y}] \mid x_{\mathbf{a}}(\mathbf{i})=\mathbf{i}\right\},  \tag{16.20b}\\
\mathbf{i}=\mathbf{i}_{n}=\left(i 1_{n}, \ldots, i 1_{n}\right) . \tag{16.21}
\end{gather*}
$$

Notice that $D[\mathfrak{r}, \mathfrak{y}]=G_{\mathbf{A}} \cap C[\mathfrak{r x}, \mathfrak{r y}]$ with $C[$,$] of type (1.17).$
We now take a fractional ideal $\mathfrak{b}$ in $F$, and recall that

$$
\begin{equation*}
G_{\mathbf{A}}=P_{\mathbf{A}} D_{0}\left[\mathfrak{b}^{-1}, \mathfrak{b}\right] . \tag{16.22}
\end{equation*}
$$

We already stated the equality $G_{\mathbf{A}}=P_{\mathbf{A}} D\left[\mathfrak{b}^{-1}, \mathfrak{b}\right]$ at the end of Lemma 1.9 , but we can take $D_{0}$ in place of $D$, since the equality at any archimedean prime holds by virtue of [S97, Propositions 6.13, 7.2, and 7.12]. Thus every element $x$ of $G_{\mathbf{A}}$ belongs to $p D_{0}\left[\mathfrak{b}^{-1}, \mathfrak{b}\right]$ for some $p \in P_{\mathbf{A}}$.

Now we fix a weight $k$ which may be integral or half-integral, and make the following convention: pr means the identity map of $G_{\mathbf{A}}$ onto itself if $k$ is integral; otherwise it is the map of (16.13). We are going to define various functions on $G_{\mathbf{A}}$ or $M_{\mathbf{A}}$ according as $k$ is integral or half-integral. For example, we define a real positive number $\varepsilon(x)$ and an ideal $\mathrm{il}_{\mathfrak{b}}(x)$ by

$$
\begin{align*}
& \varepsilon(x)=\left|\operatorname{det}\left(d_{p} d_{p}^{*}\right)\right|_{\mathbf{A}} \quad \text { and } \quad \mathrm{il}_{\mathfrak{b}}(x)=\operatorname{det}\left(d_{p}\right) \mathfrak{r} \quad \text { if } \quad \operatorname{pr}(x) \in p D_{0}\left[\mathfrak{b}^{-1}, \mathfrak{b}\right]  \tag{16.23}\\
& \text { with } p \in P_{\mathbf{A}}
\end{align*}
$$

where $\left|\left.\right|_{\mathbf{A}}\right.$ is the idele norm on $F_{\mathbf{A}}^{\times}$. These are well-defined, and
(16.23a) $\quad \varepsilon\left(x_{\mathbf{a}}\right)=\left|j_{x}(\mathbf{i})\right|^{2 \mathbf{a}}, \quad \varepsilon\left(x_{\mathbf{h}}\right)=N\left(\mathrm{il}_{\mathbf{b}}(x)\right)^{-2 /[K: F]} \quad\left(x \in G_{\mathbf{A}}\right)$.

For these see [S97, Lemma 18.5]. Clearly $\varepsilon(\pi x)=\varepsilon(x)$ for $\pi \in P$.
In addition to $\mathfrak{b}$ and $k$, we take an integral ideal $\mathfrak{c}$ in $F$ and a Hecke character $\chi$ of $K$ satisfying the following conditions:
(16.24a) $\quad \chi_{\mathbf{a}}(x)=x_{\mathbf{a}}^{\ell}\left|x_{\mathbf{a}}\right|^{i \kappa-\ell}$ with $\kappa \in \mathbf{R}^{\mathbf{a}}$ such that $\sum_{v \in \mathbf{a}} \kappa_{v}=0$ and

$$
\ell= \begin{cases}{[k]} & (\text { Case SP }) \\ \left(k_{v}-k_{v \rho}\right)_{v \in \mathbf{a}} & (\text { Case UT })\end{cases}
$$

$$
\chi_{v}(a)=1 \text { if } v \in \mathbf{h}, a \in \mathfrak{r}_{v}^{\times}, \quad \text { and } a-1 \in \mathfrak{r}_{v} \mathfrak{c}_{v}
$$

(16.24c) $D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right] \subset D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]$ if $k$ is half-integral.

Thus $\mathfrak{c} \subset 4 \mathfrak{g}$ if $k$ is half-integral. Hereafter until the end of this section we put
(16.25a) $\widetilde{D}=D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right] \quad$ and $\quad \widetilde{D}_{0}=D_{0}\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]$ if $k \in \mathbf{Z}^{\mathbf{b}}$,
(16.25b) $\quad \widetilde{D}=\left\{x \in M_{\mathbf{A}} \mid \operatorname{pr}(x) \in D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]\right\} \quad$ and

$$
\widetilde{D}_{0}=\left\{x \in \widetilde{D} \mid \operatorname{pr}(x) \in D_{0}\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]\right\} \quad \text { if } k \notin \mathbf{Z}^{\mathbf{b}}
$$

We view $P_{\mathbf{A}}$ as a subgroup of $M_{\mathbf{A}}$ by the map $r_{P}$ of (16.15). Then $P_{\mathbf{A}} \widetilde{D}$ is a subset of $M_{\mathbf{A}}$ if $k \notin \mathbf{Z}^{\mathbf{b}}$. Notice that $P_{\mathbf{A}} \widetilde{D}_{0}=P_{\mathbf{A}} \widetilde{D} \subset \mathfrak{M}$, and so $j_{\alpha}^{k}$ is meaningful for $\alpha \in P_{\mathbf{A}} \widetilde{D}$.

Next we define a function $\mu$ on $G_{\mathbf{A}}$ or $M_{\mathbf{A}}$ as follows:

$$
\begin{align*}
& \mu(x)=0 \quad \text { if } x \notin P_{\mathbf{A}} \widetilde{D},  \tag{16.26a}\\
& \mu(x)=\chi_{\mathbf{h}}\left(\operatorname{det}\left(d_{p}\right)\right)^{-1} \chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{w}\right)\right)^{-1} j_{x}^{k}(\mathbf{i})^{-1}\left|j_{x}(\mathbf{i})\right|^{m-i \kappa} \\
& \quad \text { if } x=p w \text { with } p \in P_{\mathbf{A}} \text { and } w \in \widetilde{D},
\end{align*}
$$

where $m=k$ in Case SP and $m=\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$ in Case UT. Our Eisenstein series $E_{\mathbf{A}}(x, s)$ is defined for $(x, s) \in G_{\mathbf{A}} \times \mathbf{C}$ or $(x, s) \in M_{\mathbf{A}} \times \mathbf{C}$ by

$$
\begin{equation*}
E_{\mathbf{A}}(x, s)=E_{\mathbf{A}}(x, s ; \chi, \widetilde{D})=\sum_{\alpha \in A} \mu(\alpha x) \varepsilon(\alpha x)^{-s}, \quad A=P \backslash G . \tag{16.27}
\end{equation*}
$$

This is formally well-defined, since we can easily verify, employing (16.19), (16.16c), and (16.24a), that $\mu(\pi x)=\mu(x)$ for every $\pi \in P$. We investigated this series for integral $k$ in [S97, Sections 18, 19]. Our principal aim of this section is to treat the case of half-integral $k$.
16.6. Define an element $\zeta$ of $S p(n, F)_{\mathbf{A}}$ by

$$
\zeta_{\mathbf{a}}=1, \quad \zeta_{\mathbf{h}}=\left[\begin{array}{cc}
0 & -\delta^{-1} 1_{n}  \tag{16.28}\\
\delta 1_{n} & 0
\end{array}\right]
$$

with an element $\delta$ of $F_{\mathbf{h}}^{\times}$such that $\delta \mathfrak{g}=\mathfrak{d}$. From (6.29) we easily see that

$$
\begin{equation*}
\zeta D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]=D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right] \zeta, \quad \zeta D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right] \cup D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right] \subset C^{\theta} . \tag{16.29}
\end{equation*}
$$

We then define an element $\widetilde{\zeta}$ of $M_{\mathrm{A}}$ by

$$
\begin{equation*}
\operatorname{pr}(\tilde{\zeta})=\zeta \quad \text { and } \quad h(\tilde{\zeta}, z)=1 \tag{16.30}
\end{equation*}
$$

Notice that the condition $\operatorname{pr}(\widetilde{\zeta})=\zeta$ implies that $\widetilde{\zeta} \in \mathfrak{M}$, and so $h(\widetilde{\zeta}, z)$ is meaningful; then in view of (16.16a), the condition $h(\widetilde{\zeta}, z)=1$ determines $\widetilde{\zeta}$ uniquely.

Now let $C$ be an open subgroup of $D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]$, and $\varphi$ a function on $G_{\mathbf{A}}$ or $M_{\mathbf{A}}$ such that $\varphi(\alpha x w)=j_{w}^{k}(\mathbf{i})^{-1} \varphi(x)$ if $\alpha \in G, \operatorname{pr}(w) \in C$, and $w(\mathbf{i})=\mathbf{i}$. Define a function $g$ on $\mathcal{H}$ by $g(x(\mathbf{i}))=j_{x}^{k}(\mathbf{i}) \varphi(x)$ for every $x$ such that $\operatorname{pr}(x) \in C$. This is well-defined independently of $C$ and $g \|_{k} \gamma=g$ for every $\gamma \in G \cap C$. In Case SP we can show that the correspondence $\varphi \mapsto g$ gives a bijection from the set of all such $\varphi$ 's (with a fixed $C$ ) onto the set of all functions $g$ on $\mathcal{H}$ such that $g \|_{k} \gamma=g$ for every $\gamma \in G \cap C$. In Case UT, in order to obtain a bijection, we have to associate several functions on $\mathcal{H}$ to $\varphi$; for details, see [S97, Lemma 10.8]; see also $\S 20.1$ below.

Assuming $k$ to be integral, define $\varphi^{\prime}$ on $G_{\mathbf{A}}$ by $\varphi^{\prime}(x)=\varphi\left(x \omega^{-1}\right)$ for $x \in G_{\mathbf{A}}$ with $\omega \in\left(G_{1}\right)_{\mathbf{h}}$. By strong approximation in $G_{1}$, we can find an element $\alpha \in$ $G \cap C \omega$. Define $g^{\prime}$ on $\mathcal{H}$ corresponding to $\varphi^{\prime}$ by the above principle. Then we can easily verify that

$$
\begin{equation*}
g^{\prime}=g \|_{k} \alpha \tag{16.31a}
\end{equation*}
$$

Next, assuming $k$ to be half-integral, by strong approximation we can find an element $\zeta_{0}$ of $G$ such that $\zeta_{0} \in G \cap C \zeta^{-1}$. Clearly $\zeta_{0} \in C^{\theta}$. Define $\varphi^{\prime}$ on $M_{\mathbf{A}}$ by $\varphi^{\prime}(x)=\varphi(x \widetilde{\zeta})$ for $x \in M_{\mathbf{A}}$ and define $g^{\prime}$ on $\mathcal{H}$ corresponding to $\varphi^{\prime}$ by the above principle. Then we have

$$
\begin{equation*}
g^{\prime}=g \|_{k} \zeta_{0} \tag{16.31b}
\end{equation*}
$$

To show this, take $x \in M_{\mathbf{A}}$ so that $\operatorname{pr}(x) \in C \cap \zeta C \zeta^{-1}$; put $y=\widetilde{\zeta}^{-1} x \widetilde{\zeta}$. Then $\left(g^{\prime} \|_{k} x\right)(\mathbf{i})=\varphi^{\prime}(x)=\varphi(x \widetilde{\zeta})=\varphi(\widetilde{\zeta} y)=\varphi\left(\zeta_{0} \widetilde{\zeta} y\right)=\left(g \|_{k} \zeta_{0} \widetilde{\zeta} y\right)(\mathbf{i})$ since $\operatorname{pr}\left(\widetilde{\zeta_{0}} \widetilde{\zeta} y\right) \in C$. Now $j^{k}\left(\zeta_{0} \widetilde{\zeta} y, z\right)=j^{k}\left(\zeta_{0} x \widetilde{\zeta}, z\right)=j^{k}\left(\zeta_{0} x, z\right) j^{k}(\widetilde{\zeta}, z)=j^{k}\left(\zeta_{0} x, z\right)$, since $j^{k}(\widetilde{\zeta}, z)=1$. Then $\left(g \|_{k} \zeta_{0} \widetilde{\zeta} y\right)(\mathbf{i})=\left(g \|_{k} \zeta_{0} x\right)(\mathbf{i})$, which proves (16.31b).
16.7. Returning to $E_{\mathbf{A}}$ of (16.27), we note, in both Cases SP and UT, that

$$
\begin{equation*}
E_{\mathbf{A}}(\alpha x w, s)=\chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{w}\right)\right)^{-1} j_{w}^{k}(\mathbf{i})^{-1} E_{\mathbf{A}}(x, s) \text { if } \alpha \in G \text { and } w \in \widetilde{D}_{0} \tag{16.32}
\end{equation*}
$$

In Case SP we now define a function $E_{\mathbf{A}}^{*}(x, s)$ on $G_{\mathbf{A}} \times \mathbf{C}$ or $M_{\mathbf{A}} \times \mathbf{C}$ by

$$
E_{\mathbf{A}}^{*}(x, s)=\chi(\delta)^{-n} \cdot \begin{cases}E_{\mathbf{A}}(x \zeta, s) & \left(x \in G_{\mathbf{A}}, k \in \mathbf{Z}^{\mathbf{a}}\right)  \tag{16.33}\\ E_{\mathbf{A}}(x \widetilde{\zeta}, s) & \left(x \in M_{\mathbf{A}}, k \notin \mathbf{Z}^{\mathbf{a}}\right) .\end{cases}
$$

This is independent of the choice of $\delta$ in (16.28). We are going to study the Fourier expansion of $E_{\mathbf{A}}^{*}(x, s)$. This was essentially done for integral $k$ in [S97, Section 18] (see Remark 16.12 below), and so we consider here only Case SP, putting our emphasis on the case of half-integral $k$. First, by the principle of $\S 16.7$ we define functions $E(z, s)$ and $E^{*}(z, s)$ of $(z, s) \in \mathcal{H} \times \mathbf{C}$ so that

$$
\begin{equation*}
E(x(\mathbf{i}), s)=j_{x}^{k}(\mathbf{i}) E_{\mathbf{A}}(x, s), \quad E^{*}(x(\mathbf{i}), s)=j_{x}^{k}(\mathbf{i}) E_{\mathbf{A}}^{*}(x, s) \text { if } \operatorname{pr}(x) \in G_{\mathbf{a}} \tag{16.34}
\end{equation*}
$$

We consider $E(z, s)$ in both Cases SP and UT; as for $E^{*}$, we define it, for the moment, only in Case SP. By (16.31a, b) we have

$$
\begin{equation*}
E^{*}(z, s)=\chi(\delta)^{-n} j_{\zeta_{0}}^{k}(z)^{-1} E\left(\zeta_{0} z, s\right) \tag{16.35}
\end{equation*}
$$

Now we have

$$
\begin{equation*}
E(z, s)=\sum_{\alpha \in A_{0}} N\left(\mathrm{il}_{\mathfrak{b}}(\alpha)\right)^{u s} \chi[\alpha] \delta(z)^{s \mathbf{a}-(m-i \kappa) / 2} \|_{k} \alpha, \quad A_{0}=P \backslash\left(G \cap P_{\mathbf{A}} \widetilde{D}\right) \tag{16.36}
\end{equation*}
$$

Here $m$ is as in (16.26b); $u=2$ in Case SP and $u=1$ in Case UT; $\delta(z)=$ $\left(\operatorname{det}\left((i / 2)\left(z^{*}-z\right)\right)_{v}\right)_{v \in \mathbf{a}}$, which is consistent with (3.21) and (13.5); and $\chi[\alpha]$ is an element of $\mathbf{T}$ defined for $\alpha \in G \cap P_{\mathbf{A}} \widetilde{D}$ by

$$
\chi[\alpha]= \begin{cases}\chi_{\mathbf{a}}\left(\operatorname{det}\left(d_{\alpha}\right)\right) \chi^{*}\left(\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{b}}(\alpha)^{-1}\right) & \text { if } \mathfrak{c} \neq \mathfrak{g}  \tag{16.37}\\ \chi^{*}\left(\mathrm{il}_{\mathfrak{b}}(\alpha)\right)^{-1} & \text { if } \mathfrak{c}=\mathfrak{g}\end{cases}
$$

(There are two $\delta$ 's. However, $\delta$ in (16.35) never appears together with $z$, and there will be no fear of confusion.) Notice that by Lemma $1.11(3), \operatorname{det}\left(d_{\alpha}\right) \neq 0$ and $\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{b}}(\alpha)^{-1}$ is prime to $\mathfrak{c}$ if $\mathfrak{c} \neq \mathfrak{g}$. To prove (16.36), take $x$ in $G_{\mathbf{A}}$ or $M_{\mathbf{A}}$ so that $\operatorname{pr}(x) \in G_{\mathbf{a}}$ and $x(\mathbf{i})=z$; let $\alpha \in G$. We are going to calculate each term of (16.27). By (16.26a), $\mu(\alpha x) \neq 0$ only if $\alpha \in P_{\mathbf{A}} \widetilde{D}$. Thus $A$ of (16.27) can be replaced by $A_{0}$ of (16.36). Put $\alpha x=p w$ with $p \in P_{\mathbf{A}}$ and $w \in \widetilde{D}$. By (16.23a), $\varepsilon(\alpha x)=\varepsilon\left((\alpha x)_{\mathbf{h}}\right) \varepsilon\left((\alpha x)_{\mathbf{a}}\right), \varepsilon\left((\alpha x)_{\mathbf{h}}\right)=N\left(\mathrm{il}_{\mathfrak{b}}(\alpha)\right)^{-u}$, and $\varepsilon\left((\alpha x)_{\mathbf{a}}\right)=\left|j_{\alpha x}(\mathbf{i})\right|^{2 \mathbf{a}}$. Thus, by (16.26b),

$$
\begin{align*}
& \mu(\alpha x) \varepsilon(\alpha x)^{-s}  \tag{16.37a}\\
= & N\left(\mathrm{il}_{\mathfrak{b}}(\alpha)\right)^{u s} \chi_{\mathbf{h}}\left(\operatorname{det}\left(d_{p}\right)\right)^{-1} \chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{w}\right)\right)^{-1}\left|j_{\alpha x}(\mathbf{i})\right|^{m-i \kappa-2 s \mathbf{a}} j_{\alpha x}^{k}(\mathbf{i})^{-1} .
\end{align*}
$$

Assuming $\mathfrak{c} \neq \mathfrak{g}$, we see that $d_{\alpha}$ is invertible and $\left(d_{\alpha}\right)_{\mathbf{h}}=\left(d_{p} d_{w}\right)_{\mathbf{h}}$, and hence $\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{b}}(\alpha)^{-1}=\operatorname{det}\left(d_{\alpha}\right) \operatorname{det}\left(d_{p}\right)^{-1} \mathfrak{g}=\operatorname{det}\left(d_{w}\right) \mathfrak{g}$. Therefore $\chi^{*}\left(\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{b}}(\alpha)^{-1}\right)$ $=\left(\chi_{\mathbf{h}} / \chi_{\mathrm{c}}\right)\left(\operatorname{det}\left(d_{w}\right)\right)$, so that

$$
\chi[\alpha]=\chi_{\mathbf{h}}\left(\operatorname{det}\left(d_{p} d_{w}\right)\right)^{-1}\left(\chi_{\mathbf{h}} / \chi_{\mathfrak{c}}\right)\left(\operatorname{det}\left(d_{w}\right)\right)=\chi_{\mathbf{h}}\left(\operatorname{det}\left(d_{p}\right)\right)^{-1} \chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{w}\right)\right)^{-1} .
$$

Now $j_{\alpha}^{k}(z) j_{x}^{k}(\mathbf{i})=j_{\alpha x}^{k}(\mathbf{i})$ and $\delta(\alpha z)=\left|j_{\alpha x}(\mathbf{i})\right|^{-2}$. Combining all these and (16.34), we obtain (16.36) when $\mathfrak{c} \neq \mathfrak{g}$. The case $\mathfrak{c}=\mathfrak{g}$ can be handled in the same manner.

Next put $\Gamma=G \cap D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]$. By [S97, Lemma 18.7 (4)] there exists a subset $B$ of $G \cap\left(\prod_{v \nmid \mathfrak{c}} Q_{v} \cdot D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]\right)$ such that

$$
\begin{equation*}
G \cap P_{\mathbf{A}} \widetilde{D}=\bigsqcup_{\beta \in B} P \beta \Gamma . \tag{16.38}
\end{equation*}
$$

Let $\mathcal{R}_{\beta}$ be a complete set of representatives for $\left(P \cap \beta \Gamma \beta^{-1}\right) \backslash \beta \Gamma$. Clearly $P \backslash(G \cap$ $P_{\mathbf{A}} \widetilde{D}$ ) can be given by $\bigsqcup_{\beta \in B} \mathcal{R}_{\beta}$. Since $\mathrm{il}_{\mathfrak{b}}(\alpha)=\mathrm{il}_{\mathfrak{b}}(\beta)$ if $\alpha \in \beta \Gamma$, we have

$$
\begin{equation*}
E(z, s)=\sum_{\beta \in B} N\left(\mathrm{il}_{\mathrm{b}}(\beta)\right)^{u s} \sum_{\alpha \in \mathcal{R}_{\beta}} \chi[\alpha] \delta(z)^{s \mathbf{a}-(m-i \kappa) / 2} \|_{k} \alpha . \tag{16.39}
\end{equation*}
$$

Recall that in Case SP the ideals $\mathrm{il}_{\mathfrak{b}}(\beta)$ for $\beta \in B$ form a complete set of representatives for the ideal classes of $F$ (see [S97, Lemma 18.7 (5)]).

Take $F=\mathbf{Q}$, for example; then $\kappa=0$ and $k \in 2^{-1} \mathbf{Z}$; we can take $B=\{1\}$, and $\mathfrak{c}=c \mathbf{Z}$ with $0<c \in \mathbf{Z}$. Thus, as a special case of (16.39), we have

$$
\begin{equation*}
E(z, s)=\sum_{\gamma \in(P \cap \Gamma) \backslash \Gamma} \chi_{c}\left(\operatorname{det}\left(d_{\gamma}\right)\right)^{-1} \delta(z)^{s-k / 2} \|_{k} \gamma, \tag{16.40}
\end{equation*}
$$

where we understand that $\chi_{c}\left(\operatorname{det}\left(d_{\alpha}\right)\right)=1$ if $c=1$. Thus $\chi_{c}$ may be viewed as a Dirichlet character which is even or odd according as $[k]$ is even or odd.

Now the convergence of the series of (16.27) can be reduced to that of $\sum_{\alpha \in \mathcal{R}_{\beta}}$ of (16.39), and to that of $\sum_{\alpha \in \mathcal{R}}\left|\delta(\alpha z)^{s \mathbf{a}}\right|$, where $\mathcal{R}=\left(P \cap \Gamma^{\prime}\right) \backslash \Gamma^{\prime}$ with a congruence subgroup $\Gamma^{\prime}$ of $G$. As noted in [S97, Proposition A3.7 and §A3.9], the last series is convergent for $\operatorname{Re}(s)>n$ in Case UT and $\operatorname{Re}(s)>(n+1) / 2$ in Case SP. Thus the series for $E(x, s)$ is convergent in that domain.
16.8. Since $j_{w}^{k}$ for $w \in C^{\theta}$ is a factor of automorphy in the ordinary sense, from (16.32) we can easily derive that

$$
\begin{align*}
& E_{\mathbf{A}}^{*}(\alpha x w, s)=\chi_{\mathfrak{c}}\left(\operatorname{det}\left(a_{w}\right)\right)^{-1} j_{w}^{k}(\mathbf{i})^{-1} E_{\mathbf{A}}^{*}(x, s)  \tag{16.41}\\
& \text { if } \alpha \in G \text { and } \operatorname{pr}(w) \in D_{0}\left[\mathfrak{d}^{-2} \mathfrak{b} \mathfrak{c}, \mathfrak{d}^{2} \mathfrak{b}^{-1}\right]
\end{align*}
$$

and so it has a Fourier expansion of the form

$$
E_{\mathbf{A}}^{*}\left(r_{P}\left(\begin{array}{cc}
q & \sigma \widehat{q}  \tag{16.42}\\
0 & \widehat{q}
\end{array}\right)\right)=\sum_{h \in S} c(h, q, s) \mathbf{e}_{\mathbf{A}}^{n}(h \sigma) \quad\left(q \in G L_{n}\left(F_{\mathbf{A}}\right), \sigma \in S_{\mathbf{A}}\right)
$$

with $c(h, q, s) \in \mathbf{C}$. The principle of such an expansion is stated in [S97, Proposition 18.3] for integral $k$, but the case of half-integral $k$ is similar. (Cf. also Proposition 20.2 below.) Now write an element $z$ of $\mathfrak{H}^{\mathbf{a}}$ in the form $z=x+i y$ with $x, y \in S_{\mathrm{a}}$ with $y_{v}>0$ for every $v \in \mathbf{a}$. Take $q$ and $\sigma$ in (16.42) so that $q_{\mathrm{h}}=1, q_{\mathrm{a}}=y^{1 / 2}, \sigma_{\mathrm{h}}=0$, and $\sigma_{\mathrm{a}}=x$. Write simply $y^{1 / 2}$ for such a $q$. Then from (16.34) we easily obtain

$$
\begin{equation*}
E^{*}(x+i y, s)=\operatorname{det}(y)^{-k / 2} \sum_{h \in S} c\left(h, y^{1 / 2}, s\right) \mathbf{e}_{\mathbf{a}}^{n}(h x) . \tag{16.43}
\end{equation*}
$$

To obtain the explicit form of $c(h, q, s)$, we first put

$$
\begin{gather*}
\xi\left(g, h ; s, s^{\prime}\right)=\int_{S_{v}} \mathbf{e}_{v}^{n}(-h x) \operatorname{det}(x+i g)^{-s} \operatorname{det}(x-i g)^{-s^{\prime}} d x  \tag{16.44}\\
\left(s, s^{\prime} \in \mathbf{C} ; 0<g \in S_{v}, h \in S_{v}, v \in \mathbf{a}\right) \\
\Xi\left(y, w ; t, t^{\prime}\right)=\prod_{v \in \mathbf{a}} \xi\left(y_{v}, w_{v} ; t_{v}, t_{v}^{\prime}\right) \quad\left(t, t^{\prime} \in \mathbf{C}^{\mathbf{a}}, y \in S_{\mathbf{a}}, y_{v}>0, w \in S_{\mathbf{a}}\right) .
\end{gather*}
$$

Here for $s \in \mathbf{C}$ and $z \in \mathfrak{H}_{n}$ we choose the branches of $\operatorname{det}(z)^{s}$ and $\overline{\operatorname{det}(z)}^{s}$ so that their values at $z=i 1_{n}$ are $i^{n s}$ and $i^{-n s}$, respectively, where $i^{\alpha}$ is defined by

$$
\begin{equation*}
i^{\alpha}=\exp (\pi i \alpha / 2) \quad(\alpha \in \mathbf{C}) \tag{16.45}
\end{equation*}
$$

The function $\xi$ was investigated in [S82]. We note here only that the integral of (16.44) is convergent for sufficiently large $\operatorname{Re}\left(s+s^{\prime}\right)$, and can be continued as a meromorphic function of $\left(s, s^{\prime}\right)$ to the whole $\mathbf{C}^{2}$. We also put, for $\tau \in \widetilde{S}$,

$$
\begin{aligned}
& \alpha_{\mathfrak{c}}^{0}(\tau, s, \chi)=\prod_{v \nmid c} \sum_{\sigma \in S_{v} / S(\mathfrak{g})_{v}} \mathbf{e}_{v}^{n}\left(-\delta_{v}^{-1} \tau \sigma\right) \chi^{*}\left(\nu_{0}(\sigma)\right) \nu(\sigma)^{-s}, \\
& \alpha_{\mathfrak{a}}^{1}(\tau, s, \chi)=\prod_{v \nmid c} \sum_{\sigma \in S_{v} / S(\mathfrak{g})_{v}} \mathbf{e}_{v}^{n}\left(-\delta_{v}^{-1} \tau \sigma\right) \chi^{*}\left(\nu_{0}(\sigma)\right) \omega\left(\delta_{v}^{-1} \sigma\right) \nu(\sigma)^{-s} .
\end{aligned}
$$

16.9. Proposition (Case SP). Suppose that $\mathfrak{c} \neq \mathfrak{g}$ and $\operatorname{det}\left(q_{v}\right)>0$ for every $v \in \mathbf{a}$; let $y={ }^{t} q_{\mathbf{a}} q_{\mathbf{a}}$. Then $c(h, q, s) \neq 0$ only if $\left({ }^{t} q h q\right)_{v} \in\left(\mathfrak{d b}^{-1} \mathfrak{c}^{-1}\right)_{v} \widetilde{S}_{v}$ with $\widetilde{S}_{v}$ of (16.1c) for every $v \in \mathbf{h}$, in which case

$$
\begin{aligned}
c(h, q, s)= & C \cdot \chi_{\mathbf{h}}(\operatorname{det}(-q))^{-1}\left|\operatorname{det}(q)_{\mathbf{h}}\right|_{\mathbf{A}}^{n+1-2 s}\left|D_{F}\right|^{-2 n s+3 n(n+1) / 4} N(\mathfrak{b} \mathbf{c})^{-n(n+1) / 2} \\
& \cdot \operatorname{det}(y)^{s \mathbf{a}+i \kappa / 2} \Xi(y, h ; s \mathbf{a}+(k+i \kappa) / 2, s \mathbf{a}-(k-i \kappa) / 2) \\
& \cdot \alpha_{\mathfrak{c}}^{e}\left(\varepsilon_{b}^{-1} \cdot{ }^{t} q h q, 2 s, \chi\right),
\end{aligned}
$$

where $C=1$ and $e=0$ if $k \in \mathbf{Z}^{\mathbf{a}}$, and $C=\mathbf{e}(n[F: \mathbf{Q}] / 8)$ and $e=1$ if $k \notin \mathbf{Z}^{\mathbf{a}} ; \varepsilon_{b}$ is an element of $F_{\mathbf{h}}^{\times}$such that $\varepsilon_{b} \mathfrak{g}=\mathfrak{b}^{-1} \mathfrak{d}$ if $k \in \mathbf{Z}^{\mathbf{a}}$, and $\varepsilon_{b}=1$ if $k \notin \mathbf{Z}^{\mathbf{a}} ; D_{F}$ is the discriminant of $F$.

If $k \notin \mathbf{Z}^{\mathbf{a}}$, our assumption (16.24c) implies that $\mathfrak{b}_{v}=\mathfrak{d}_{v}$ for $v \nmid c$. Therefore $\left(\varepsilon_{b}^{-1} \cdot{ }^{t} q h q\right)_{v} \in \widetilde{S}_{v}$ if $v \nmid c$ for both integral and half-integral $k$. The proof will be given in §A2.13.
16.10. Proposition (Case SP). With $h, q$ such that $c(h, q, s) \neq 0$ and $\varepsilon_{b}$ as in Proposition 16.9, put $r=\operatorname{rank}(h)$ and ${ }^{t} g h g=\operatorname{diag}\left[h^{\prime}, 0\right]$ with $g \in G L_{n}(F)$ and $h^{\prime} \in S^{r}$. Let $\rho_{h}$ be the Hecke character corresponding to $F\left(c^{1 / 2}\right) / F$, where $c=(-1)^{[r / 2]} \operatorname{det}\left(2 h^{\prime}\right)$, if $r>0$; let $\rho_{h}=1$ if $r=0$. Then

$$
\begin{equation*}
\alpha_{\mathfrak{c}}^{e}\left(\varepsilon_{b}^{-1} \cdot{ }^{t} q h q, 2 s, \chi\right)=\Lambda_{\mathfrak{c}}(s)^{-1} \Lambda_{h}(s) \prod_{v \in \mathbf{c}} f_{h, q, v}\left(\chi\left(\pi_{v}\right)\left|\pi_{v}\right|^{2 s+e / 2}\right) \tag{16.46}
\end{equation*}
$$

with a finite subset $\mathbf{c}$ of $\mathbf{h}$, polynomials $f_{h, q, v}$ with coefficients in $\mathbf{Z}$ independent of $\chi$, and functions $\Lambda_{c}$ and $\Lambda_{h}$ given as follows:

$$
\begin{aligned}
& \Lambda_{\mathfrak{c}}(s)=\Lambda_{\mathfrak{c}}^{n}(s)= \begin{cases}L_{\mathfrak{c}}(2 s, \chi) \prod_{i=1}^{[n / 2]} L_{\mathfrak{c}}\left(4 s-2 i, \chi^{2}\right) & \text { if } k \in \mathbf{Z}^{\mathbf{a}}, \\
\prod_{i=1}^{[(n+1) / 2]} L_{\mathfrak{c}}\left(4 s-2 i+1, \chi^{2}\right) & \text { if } k \notin \mathbf{Z}^{\mathbf{a}},\end{cases} \\
& \Lambda_{h}(s)=L_{\mathfrak{c}}\left(2 s-n+r / 2, \chi \rho_{h}\right) \prod_{i=1}^{\lfloor(n-r) / 2]} L_{\mathfrak{c}}\left(4 s-2 n+r+2 i-1, \chi^{2}\right) \\
& \text { if } 2 k_{v}+r \in 2 \mathbf{Z} \text {, } \\
& \Lambda_{h}(s)=\prod_{i=1}^{[(n-r+1) / 2]} L_{\mathfrak{c}}\left(4 s-2 n+r+2 i-2, \chi^{2}\right) \\
& \text { if } \quad 2 k_{v}+r \notin 2 \mathbf{Z} .
\end{aligned}
$$

The set $\mathbf{c}$ is determined as follows: $\mathbf{c}=\varnothing$ if $r=0$. If $r>0$, take $a \in \prod_{v \nmid c} G L_{n}\left(\mathfrak{g}_{v}\right)$ so that $\left(\varepsilon_{b}^{-1} \cdot{ }^{t} a \cdot{ }^{t} q h q a\right)_{v}=\operatorname{diag}\left[\tau_{v}, 0\right]$ with $\tau_{v} \in \widetilde{S}_{v}^{r}$ for every $v \nmid c$. (Such an $a$ is guaranteed by [S97, Lemma 13.3].) Then c consists of all the primes $v$ not dividing $\mathfrak{c}$ such that $\tau_{v}$ is not regular in the sense of §16.1.

This follows immediately from Theorem 16.2.
Next, to state a theorem concerning the analytic nature of $E(z, s)$, we first put, for $\iota=1$ or 2 , and $s \in \mathbf{C}$,

$$
\begin{equation*}
\Gamma_{0}^{\iota}(s)=1, \quad \Gamma_{n}^{\iota}(s)=\pi^{\iota n(n-1) / 4} \prod_{\nu=0}^{n-1} \Gamma(s-(\iota \nu / 2)) \quad(n>0) . \tag{16.47}
\end{equation*}
$$

16.11. Theorem (Case SP). Write simply $\Gamma_{n}$ for $\Gamma_{n}^{1}$; define $\gamma(s, h)$ and $\mathcal{G}(s)$ as follows:

$$
\begin{aligned}
\mathcal{G}(s)=\mathcal{G}_{k, \kappa}^{n}(s)=\prod_{v \in \mathbf{a}} \gamma\left(s+i \kappa_{v} / 2,\left|k_{v}\right|\right), \\
\gamma(s, h)= \begin{cases}\Gamma\left(s+\frac{h}{2}-\left[\frac{2 h+n}{4}\right]\right) \Gamma_{n}\left(s+\frac{h}{2}\right) & (n / 2 \leq h \in \mathbf{Z}, n \text { even }) \\
\Gamma_{n}\left(s+\frac{h}{2}\right) & (n / 2<h \in \mathbf{Z}, n \text { odd }), \\
\Gamma_{2 h+1}\left(s+\frac{h}{2}\right) \prod_{i=h+1}^{[n / 2]} \Gamma(2 s-i) & (0 \leq h<n / 2, h \in \mathbf{Z})\end{cases}
\end{aligned}
$$

$$
\gamma(s, h)= \begin{cases}\Gamma\left(s+\frac{h-1}{2}-\left[\frac{2 h+n-2}{4}\right]\right) \Gamma_{n}\left(s+\frac{h}{2}\right) & (n / 2<h \notin \mathbf{Z}, n \text { odd }), \\ \Gamma_{n}\left(s+\frac{h}{2}\right) & (n / 2<h \notin \mathbf{Z}, n \text { even }), \\ \Gamma_{2 h+1}\left(s+\frac{h}{2}\right) \prod_{i=[h]+1}^{[(n-1) / 2]} \Gamma\left(2 s-\frac{1}{2}-i\right) & (0<h \leq n / 2, h \notin \mathbf{Z})\end{cases}
$$

Put $\mathcal{P}(s)=\mathcal{G}(s) \Lambda_{\mathfrak{c}}^{n}(s) E(z, s)$ with $E$ of (16.36). Then $\mathcal{P}(s)$ is a meromorphic function in $s$ on the whole $\mathbf{C}$ with only finitely many poles, and each pole is simple. In particular, $\mathcal{P}$ is an entire function of $s$ if $\chi^{2} \neq 1$. If $\chi^{2}=1$, the poles of $\mathcal{P}$ are determined as follows:
(I) $\chi^{2}=1$ and $\mathfrak{c} \neq \mathfrak{g}$ : Let $m=\operatorname{Max}_{v \in \mathbf{a}}\left|k_{v}\right|$. If $m>n / 2, \mathcal{P}$ has no pole except for a possible pole at $s=(n+2) / 4$ which occurs only if $2\left|k_{v}\right|-n \in 4 \mathbf{Z}$ for every $v$ such that $2\left|k_{v}\right|>n$. If $m \leq n / 2, \mathcal{P}$ has possible poles only in the set

$$
\begin{align*}
& \{j / 2 \mid j \in \mathbf{Z},[(n+3) / 2] \leq j \leq n+1-m\} \quad \text { if } \quad k \in \mathbf{Z}^{\mathbf{a}},  \tag{i}\\
& \{(2 j+1) / 4 \mid j \in \mathbf{Z}, 1+[n / 2] \leq j \leq n+(1 / 2)-m\} \quad \text { if } \quad k \notin \mathbf{Z}^{\mathbf{a}} . \tag{ii}
\end{align*}
$$

(II) $\chi^{2}=1, \mathfrak{c}=\mathfrak{g}$, and $k \in \mathbf{Z}^{\mathbf{a}}$ : In this case each pole belongs to the set of poles described in (I) or to

$$
\begin{equation*}
\{j / 2 \mid j \in \mathbf{Z}, 0 \leq j \leq[n / 2]\} \tag{iii}
\end{equation*}
$$

where $j=0$ is unnecessary if $\chi \neq 1$.
We can derive this from Propositions 16.9 and 16.10 by means of the principle of [S97, Proposition 19.1] and the procedure described in [S97, §§19.4~19.6] in Case UT. In fact, the detailed discussion was given in [S94a, pp.565-571].
16.12. Remark. (I) In [S97, §18.6], we considered $E(x, s)$ for integral $k$ in Cases SP and UT. In Case UT we have $j_{x}^{k}=\operatorname{det}(x)^{\nu} j_{x}^{m}$ with $\nu=\left(-k_{v \rho}\right)_{v \in \mathbf{a}}$ and $m=\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$, and hence the present $j_{x}^{k}$ can be written $j_{x}^{m, \nu}$ in the notation of [S97, (10.4.3)]. Thus our $E(x, s)$ coincides with $E$ of $[S 97,(18.6 .1)]$ if we take $j_{x}^{k, \nu}$ there to be the present $j_{x}^{m, \nu}$. Also we considered there $E^{*}(x, s)=E_{\mathbf{A}}\left(x \eta_{\mathbf{h}}^{-1}, s\right)$ in both Cases SP and UT, instead of (16.33), and obtained its Fourier expansion in [S97, Proposition 18.14]. The function is different from the present $E_{\mathbf{A}}^{*}$, but the nature is the same. At any rate, both functions can be handled by the same methods.
(II) We can define $E_{\mathbf{A}}(x, s)$ also on $\left(G_{1}\right)_{\mathbf{A}}$ with $G_{1}=S U\left(\eta_{n}\right)$ in Case UT and the corresponding $E(z, s)$. This was done in [S97, Section 18]. In this case we take $k \in \mathbf{Z}^{\mathbf{a}}$. By Lemma $1.3(2), \operatorname{det}\left(d_{\alpha}\right) \in F$ for every $\alpha \in G_{1}$. Therefore we put $\mathrm{il}_{\mathfrak{b}}(\alpha)=\operatorname{det}\left(d_{p}\right) \mathfrak{g}$ in (16.23), and take $\chi$ to be a Hecke character of $F$ satisfying $\chi_{\mathbf{a}}(x)=x_{\mathbf{a}}^{k}\left|x_{\mathbf{a}}\right|^{i \kappa-k}$. To obtain the formula corresponding to (16.39) in this case, we take $G_{1}$ and $P_{1}$ in place of $G$ and $P$, and take $B$ so that $G_{1} \cap\left(P_{1}\right)_{\mathbf{A}} D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]_{1}=$ $\bigsqcup_{\beta \in B} P_{1} \beta \Gamma_{1}$, where $X_{1}=\left(G_{1}\right)_{\mathbf{A}} \cap X$ for any subgroup $X$ of $G_{\mathbf{A}}$. With this choice of $B$ the ideals $\mathrm{il}_{\mathfrak{b}}(\beta)$ for $\beta \in B$ form a complete set of repressentatives for the ideal classes of $F$; see [S97, Lemma 18.7 (5)]. Then with $\mathcal{R}_{\beta}=\left(P_{1} \cap \beta \Gamma_{1} \beta^{-1}\right) \backslash \beta \Gamma_{1}$ we have

$$
\begin{equation*}
E(z, s)=\sum_{\beta \in B} N\left(\mathrm{il}_{\mathfrak{b}}(\beta)\right)^{2 s} \sum_{\alpha \in \mathcal{R}_{\beta}} \chi[\alpha] \delta(z)^{s \mathrm{a}-(k-i \kappa) / 2} \|_{k} \alpha \tag{16.48}
\end{equation*}
$$

(III) Let $L=\sum_{i=1}^{n}\left(\mathfrak{r} e_{i}+\mathfrak{b}^{-1} e_{n+i}\right)$ with the standard basis $\left\{e_{i}\right\}_{i=1}^{2 n}$ of $K_{2 n}^{1}$ in Case UT. Take $\zeta \in K^{\times}$such that $\zeta^{\rho}=-\zeta$. Then we can easily verify that $L$ is a maximal lattice with respect to the hermitian form $\zeta \eta$ in the sense of [S97, §4.7]. Clearly $D\left[\mathfrak{b}^{-1}, \mathfrak{b}\right]=\left\{\alpha \in G_{\mathbf{A}} \mid L \alpha=L\right\}$. Therefore the results of [S97, Section 5] are applicable to $L_{v}$ and $D\left[\mathfrak{b}^{-1}, \mathfrak{b}\right] \cap G_{v}$ for every $v \in \mathbf{h}$.

## 17. Arithmeticity and near holomorphy of Eisenstein series

17.1. The purpose of this section is to study the nature of $E(z, s)$ at suitably chosen points $s_{0}$ belonging to an interval. We are going to show that the value or the residue of $E$ at $s_{0}$ is nearly holomorphic, sometimes holomorphic, and arithmetic up to a power of $\pi$. Throughout this section we put

$$
\begin{equation*}
d=[F: \mathbf{Q}] . \tag{17.1}
\end{equation*}
$$

We return to the setting of $\S \S 16.5$ and 16.7 in Cases SP and UT. However, in Case UT we take $G_{1}=S U\left(\eta_{n}\right)$ instead of $G=U\left(\eta_{n}\right)$. Therefore, in this section until Theorem 17.9, we speak of Case SU instead of Case UT. We shall return to Case UT in $\S 17.10$. The notation being as in $\S 16.7$, let $\Gamma$ be a congruence subgroup of $G_{1}$ or $\mathcal{G}$ according as $k$ is integral or half-integral, where $\mathcal{G}$ is the group defined in §14.14. As we said there, we view every congruence subgroup of $\Gamma^{\theta}$ as a congruence subgroup of $\mathcal{G}$. In Case SU we take $k \in \mathbf{Z}^{\mathbf{a}}$. We also take an element $\kappa$ of $\mathbf{R}^{\mathbf{a}}$ such that $\sum_{v \in \mathbf{a}} \kappa_{v}=0$. For $\xi=(\alpha, p) \in \mathcal{G}$ we put $a_{\xi}=a_{\alpha}, b_{\xi}=b_{\alpha}, c_{\xi}=c_{\alpha}$, and $d_{\xi}=d_{\alpha}$. Then we put

$$
\begin{gather*}
\Gamma^{P}=\left\{\gamma \in \Gamma \mid c_{\gamma}=0\right\}  \tag{17.2}\\
E(z, s ; k, \kappa, \Gamma)=\sum_{\alpha \in \Gamma^{P} \backslash \Gamma} \delta(z)^{s \mathbf{a}-(k-i \kappa) / 2} \|_{k} \alpha \quad((z, s) \in \mathcal{H} \times \mathbf{C}) .  \tag{17.3}\\
E(z, s ; k, \Gamma)=E(z, s ; k, 0, \Gamma) . \tag{17.3a}
\end{gather*}
$$

To make the sum at least formally meaningful, we have to assume

$$
\begin{equation*}
\left|j_{\gamma}(z)\right|^{i \kappa-k} j_{\gamma}^{k}(z)=1 \quad \text { for every } \gamma \in \Gamma^{P} . \tag{17.4}
\end{equation*}
$$

The series is convergent for $\operatorname{Re}(s)>(n+1) / 2$ in Case SP and $\operatorname{Re}(s)>n$ in Case UT, as noted in [S97, Proposition A3.7 and §A3.9]. If $\Gamma_{1}$ is a congruence subgroup contained in $\Gamma$, then we can easily verify that

$$
\begin{equation*}
\left[\Gamma^{P}: \Gamma_{1}^{P}\right] E(z, s ; k, \kappa, \Gamma)=\sum_{\alpha \in \Gamma_{1} \backslash \Gamma} E\left(z, s ; k, \kappa, \Gamma_{1}\right) \|_{k} \alpha . \tag{17.5}
\end{equation*}
$$

To study the properties of $E(z, s ; k, \kappa, \Gamma)$ for an arbitrary $\Gamma$, we introduce special types of congruence subgroups as follows:

$$
\begin{align*}
& \Gamma_{0}(\mathfrak{c})=G_{1} \cap D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right],  \tag{17.6a}\\
& \Gamma(\mathfrak{c})=\left\{\alpha \in \Gamma_{0}(\mathfrak{c}) \mid a_{\gamma}-1 \prec \mathfrak{r c}, b_{\alpha} \prec \mathfrak{r b}^{-1} \mathfrak{c}\right\},  \tag{17.6b}\\
& \Gamma_{u}(\mathfrak{c})=\left\{\alpha \in \Gamma_{0}(\mathfrak{c}) \mid \operatorname{det}\left(d_{\alpha}\right) \in \mathfrak{g}^{\times} W\right\}, \\
& W=\left\{x \in F^{\times} \mid x-1 \in \mathfrak{c}_{v} \text { for every } v \mid \mathfrak{c}\right\} .
\end{align*}
$$

Here we fix $\mathfrak{b}$ and assume (16.24c). We also denote by $E(z, s ; k, \chi, \mathfrak{c})$ the series of (16.36) in Case SP and that of (16.48) in Case SU; though $E$ depends also on $\mathfrak{b}$, we suppress it, since it has no essential effect on the nature of $E$.
17.2. Lemma (Cases SP and SU). (1) $\Gamma_{u}(\mathfrak{c})=\Gamma_{u}(\mathfrak{c})^{P} \Gamma(\mathfrak{c})$.
(2) Let $X$ be the set of all Hecke characters of $F$ satisfying (16.24a, b) with a fixed $\kappa$. Let $\Gamma$ be a congruence subgroup of $G_{1}$ or $\mathcal{G}$ which satisfies (17.4) and contains $\Gamma(\mathfrak{c})$. Then $X \neq \varnothing$ and

$$
\#(X)\left[\Gamma^{P}: \Gamma(\mathfrak{c})^{P}\right] E(z, s ; k, \kappa, \Gamma)=\sum_{\chi \in X} \sum_{\alpha \in \Gamma(\mathfrak{c}) \backslash \Gamma} E(z, s ; k, \chi, \mathfrak{c}) \|_{k} \alpha
$$

(3) If $\mathfrak{a}$ is a multiple of $\mathfrak{c}$, then

$$
E(z, s ; k, \chi, \mathfrak{c})=\sum_{\xi \in \Gamma_{0}(\mathfrak{a}) \backslash \Gamma_{0}(\mathfrak{c})} \chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{\xi}\right)\right)^{-1} E(z, s ; k, \chi, \mathfrak{a}) \|_{k} \xi
$$

(4) Let $\lambda_{n}=(n+1) / 2$ in Case SP and $\lambda_{n}=n$ in Case UT. Then $E(z, s ; k, \kappa, \Gamma)$ can be continued to a meromorphic function on the whole s-plane which is holomorphic for $\operatorname{Re}(s)>\lambda_{n}$. Moreover it has a pole at $s=\lambda_{n}$ only if $k=\kappa=0$, in which case it has a simple pole at $s=\lambda_{n}$ with a positive real number as its residue.

Proof. To prove (1), take $\alpha \in \Gamma_{u}(\mathfrak{c})$. Then $\operatorname{det}\left(d_{\alpha}\right) \in e W$ with $e \in \mathfrak{g}^{\times}$. By strong approximation on $S L_{n}(K)$, we can find an element $q$ of $G L_{n}(\mathfrak{r})$ such that $q-d_{\alpha} \prec \mathfrak{c r}$ and $\operatorname{det}(q)=e$. Put

$$
\beta=\left[\begin{array}{cc}
q^{*} & -d_{\alpha}^{*} b_{\alpha} q^{-1} \\
0 & q^{-1}
\end{array}\right]
$$

Then $\beta \in \Gamma_{u}(\mathfrak{c})^{P}$ and $\beta \alpha \in \Gamma(\mathfrak{c})$, which proves (1). Next, since $\Gamma(\mathfrak{c}) \subset \Gamma$, from (16.19) and (17.4) we see that $e^{[k]}|e|^{i \kappa-[k]}=1$ for every $e \in \mathfrak{g}^{\times}$such that $e-1 \in \mathfrak{c}$. Thus $X \neq \varnothing$ by [S97, Lemma 11.14]. Take any $\chi_{0} \in X$; observe that, for a $\mathfrak{g}$-ideal $\mathfrak{a}$ prime to $\mathfrak{c}$, we have $\sum_{\chi \in X} \chi^{*}(\mathfrak{a})=\#(X) \chi_{0}(\mathfrak{a})$ if $\mathfrak{a}=a \mathfrak{g}$ with $a \in W$; otherwise the sum is 0 . Now take the sum of (16.39) or (16.48) for all $\chi \in X$. For $\alpha \in \mathcal{R}_{\beta} \subset \beta \Gamma_{0}(\mathfrak{c})$ we have $\sum_{\chi \in X} \chi[\alpha] \neq 0$ only when $\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{b}}(\alpha)^{-1}=a \mathfrak{g}$ with $a \in W$, so that $\mathrm{il}_{\mathfrak{b}}(\beta)$ is principal. Since the $\mathrm{il}_{\mathfrak{b}}(\beta)$ for all $\beta \in B$ represent the ideal classes of $F$, we can take $\beta=1$, and hence $\alpha \in \Gamma_{0}(\mathfrak{c})$ and $\operatorname{det}\left(d_{\alpha}\right) \in W \mathfrak{g}^{\times}$, that is, $\alpha \in \Gamma_{u}(\mathfrak{c})$. Also $\chi[\alpha]=\chi_{\mathbf{a}}\left(\operatorname{det}\left(d_{\alpha}\right)\right) \chi^{*}\left(\operatorname{det}\left(d_{\alpha}\right) \mathfrak{g}\right)=\chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{\alpha}\right)\right)^{-1}$. Since $\Gamma_{0}(\mathfrak{c})^{P}=\Gamma_{u}(\mathfrak{c})^{P}$, we have

$$
\begin{equation*}
\sum_{\chi \in X} E(z, s ; k, \chi, \mathfrak{c})=\#(X) \sum_{\alpha \in \mathcal{R}}\left(\chi_{0}\right)_{\mathfrak{c}}\left(\operatorname{det}\left(d_{\alpha}\right)\right)^{-1} \delta(x)^{s \mathbf{a}-(k-i \kappa) / 2} \|_{k} \alpha \tag{*}
\end{equation*}
$$

with $\mathcal{R}=\Gamma_{u}(\mathfrak{c})^{P} \backslash \Gamma_{u}(\mathfrak{c})$. By (1) we can take $\Gamma(\mathfrak{c})^{P} \backslash \Gamma(\mathfrak{c})$ as $\mathcal{R}$. Then the right-hand side of $\left(^{*}\right)$ is $\#(X) E(z, s ; k, \Gamma(\mathfrak{c}))$. Combining this with (17.5), we obtain (2). To prove (3), put $\Gamma=\Gamma_{0}(\mathfrak{c})$ and $\Gamma^{\prime}=\Gamma_{0}(\mathfrak{a})$; denote by $\widetilde{D}^{\prime}$ the group of (16.25a, b) defined with $\mathfrak{a}$ in place of $\mathfrak{c}$. By strong approximation we have $\widetilde{D} \subset \widetilde{D}^{\prime} G_{1}$, and so $\widetilde{D}=\widetilde{D}^{\prime} \Gamma$. Let $T_{\widetilde{D}}=\Gamma^{\prime} \backslash \Gamma, A_{0}=P \backslash\left(G \cap P_{\mathbf{A}} \widetilde{D}\right)$, and $A_{0}^{\prime} \underset{\sim}{=} P \backslash\left(G \cap P_{\mathbf{A}} \widetilde{D}^{\prime}\right)$. Since $P_{\mathbf{A}} \cap \widetilde{D}=P_{\mathbf{A}} \cap \widetilde{D}^{\prime}$, we easily see that $P_{\mathbf{A}} \widetilde{D}=\bigsqcup_{\xi \in T} P_{\mathbf{A}} \widetilde{D}^{\prime} \xi$. Therefore $A_{0}$ can be given by $\bigsqcup_{\xi \in T} A_{0}^{\prime} \xi$. From this and (16.36) or (16.48) we obtain (3), since $\chi[\alpha \xi]=\chi[\alpha] \chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{\xi}\right)\right)^{-1}$ for $\xi \in \Gamma$.

As for (4), the holomorphy for $\operatorname{Re}(s)>\lambda_{n}$ follows from the convergence of our series in that domain. Now $E(z, s ; k, \chi, \mathfrak{c})$ has meromorphic continuation to the whole $s$-plain by Theorem 16.11 in Case SP and by [S97, Theorem 19.7] in Case SU. Also, from these theorems we see that $E(z, s ; k, \chi, \mathfrak{c})$ is holomorphic at $s=\lambda_{n}$ except when $k=0$ and $\chi=1$, in which case it has a simple pole at $s=\lambda_{n}$ with a positive number as its residue. Combining this with the equality of (2), we obtain (4).
17.3. Hereafter we assume that $\kappa=0$, which means that a Hecke character $\chi$ of $F$ satisfying (16.24a) is of finite order. Let us now put

$$
\begin{align*}
& D(z, s ; k, \chi, \mathfrak{c})  \tag{17.7}\\
= & E(z, s ; k, \chi, \mathfrak{c}) \cdot\left\{\begin{array}{l}
L_{\mathfrak{c}}(2 s, \chi) \prod_{i=1}^{[n / 2]} L_{\mathfrak{c}}\left(4 s-2 i, \chi^{2}\right) \\
\prod_{i=1}^{[(n+1) / 2]} L_{\mathfrak{c}}\left(4 s-2 i-1, \chi^{2}\right) \quad\left(\text { Case SP, } k \in \mathbf{Z}^{\mathbf{a}}\right), \\
\prod_{i=0}^{n-1} L_{\mathfrak{c}}\left(2 s-i, \chi \theta^{i}\right) \quad\left(\text { Case SP, } k \notin \mathbf{Z}^{\mathbf{a}}\right),
\end{array}\right.
\end{align*}
$$

where $\theta$ is the quadratic Hecke character of $F$ corresponding to $K / F$.
We begin our investigation of $E$ at special values of $s$ by considering the special case in which $k=\mu \mathbf{a}$ with $\mu \in 2^{-1} \mathbf{Z}$; naturally $\mu \in \mathbf{Z}$ in Case SU . We then define $E^{*}(z, s)$ in Case SP as in $\S 16.7$, and in Case SU we put

$$
\begin{equation*}
E^{*}(z, s)=E(z, s ; \mu \mathbf{a}, \chi, \mathfrak{c}) \|_{\mu \mathbf{a}} \eta \quad(\text { Case } \mathrm{SU}) \tag{17.8}
\end{equation*}
$$

Here we assume that $\mathfrak{c} \neq \mathfrak{g}$ and $\chi_{\mathbf{a}}(x)=\operatorname{sgn}\left(x_{\mathbf{a}}\right)^{[\mu] \mathbf{a}}$. We have a Fourier expansion

$$
\begin{align*}
E^{*}(x+i y, s)= & \sum_{h \in S} c_{h}(y, s) \mathbf{e}_{\mathbf{a}}^{n}(h x)  \tag{17.9}\\
& \text { with } \quad c_{h}(y, s)=\operatorname{det}(y)^{-\mu \mathbf{a} / 2} c\left(h, y^{1 / 2}, s\right)
\end{align*}
$$

as observed in (16.43) and [S97, Lemma 18.7 (2)]. The formula for $c(h, \cdots)$ in Case SU is given in [S97, Propositions 18.14 and 19.2]; these correspond to Propositions 16.9 and 16.10.

Now, if the series $E(z, s)$ of (16.36) or (17.3) with $k=\mu \mathbf{a}$ is convergent at $s=$ $\mu / 2$, then clearly $E(z, \mu / 2)$ is holomorphic in $z$, and so it belongs to $\mathcal{M}_{\mu \mathbf{a}}$. It can happen that $E(z, s)$ is not convergent at $s=\mu / 2$, but that analytic continuation allows us to speak of $E(z, \mu / 2)$. Then we can ask the following questions:
(R1) When is $E(z, \mu / 2)$ meaningful? If $E(z, \mu / 2)$ is meaningful, is it holomorphic in $z$ ? If so, is it an element of $\mathcal{M}_{\mu \mathbf{a}}(\overline{\mathbf{Q}})$ ?
(R2) If $E(z, \mu / 2)$ is not holomorphic in $z$, can we say something about its analytic nature? What happens if we take a more general $k$ ?
(R3) Are there more values of $s$ for which we can describe the nature of $E(z, s)$ ?
We can ask similar questions by taking $D$ of (17.7) in place of $E$. We shall answer these questions in Theorems 17.7, 17.8, and 17.9. In particular, we shall see that $E(z, \mu / 2)$ can be nonholomorphic but nearly holomorphic in certain cases.
17.4. We are going to consider the behavior of $c_{h}(y, s)$ at $s=\mu / 2$. To recall some of the properties of the function $\xi$ of (16.44), we first put

$$
\begin{align*}
& \lambda=\lambda_{n}=(n+1) / 2, \quad \iota=1 \quad(\text { Case SP }),  \tag{17.10a}\\
& \lambda=\lambda_{n}=n, \quad \iota=2 \quad(\text { Case SU }) \tag{17.10b}
\end{align*}
$$

In [S82] we obtained a function $\omega(g, h ; \alpha, \beta)$ defined for $(g, h ; \alpha, \beta)$ as in (16.44), holomorphic in $(\alpha, \beta) \in \mathbf{C}^{2}$, with which we have

$$
\begin{align*}
\xi(g, h ; \alpha, \beta)= & i^{n \beta-n \alpha} 2^{\tau} \pi^{\varepsilon} \Gamma_{t}^{\iota}(\alpha+\beta-\lambda) \Gamma_{n-q}^{\iota}(\alpha)^{-1} \Gamma_{n-p}^{\iota}(\beta)^{-1}  \tag{17.11}\\
& \cdot \operatorname{det}(g)^{\lambda-\alpha-\beta} \delta_{+}(h g)^{\alpha-\lambda+\iota q / 4} \delta_{-}(h g)^{\beta-\lambda+\iota p / 4} \\
& \cdot \omega(2 \pi g, h ; \alpha, \beta) .
\end{align*}
$$

Here $p$ (resp. $q$ ) is the number of positive (resp. negative) eigenvalues of $h$ and $t=n-p-q ; \delta_{+}(x)$ is the product of all positive eigenvalues of $x$ and $\delta_{-}(x)=\delta_{+}(-x) ; \Gamma_{n}^{\iota}$ is defined by (16.47);

$$
\begin{aligned}
& \tau=(2 p-n) \alpha+(2 q-n) \beta+n+t \lambda+(\iota p q / 2) \\
& \varepsilon=p \alpha+q \beta+t+(\iota / 2)\{t(t-1)-p q\}
\end{aligned}
$$

In particular we have

$$
\begin{align*}
& \xi(g, h ; \alpha, 0)=2^{1-\lambda} i^{-n \alpha}(2 \pi)^{n \alpha} \Gamma_{n}^{\iota}(\alpha)^{-1} \operatorname{det}(h)^{\alpha-\lambda} \mathbf{e}^{n}(i g h) \quad \text { if } h>0  \tag{17.12}\\
& \xi(g, 0 ; \alpha, \beta)=i^{n \beta-n \alpha} 2^{n(\lambda+1-\alpha-\beta)} \pi^{n \lambda} \frac{\Gamma_{n}^{\iota}(\alpha+\beta-\lambda)}{\Gamma_{n}^{\iota}(\alpha) \Gamma_{n}^{\iota}(\beta)} \operatorname{det}(g)^{\lambda-\alpha-\beta}  \tag{17.13}\\
& \lim _{s \rightarrow 0} \xi(g, h ; \lambda+s, s)=2^{\sigma} i^{-n \lambda} \pi^{n \lambda} \Gamma_{n}^{\iota}(\lambda)^{-1} \mathbf{e}^{n}(i g h) \text { if } q=0  \tag{17.14}\\
& \text { with } \sigma=[(n+p) / 2] \text { in Case } S P \text { and } \sigma=p \text { in Case SU. }
\end{align*}
$$

For these see $[\mathrm{S} 82,(4.34 . \mathrm{K}),(4.35 . \mathrm{K}),(1.31)]$. To state one more formula, we take an indeterminate $T$ and define polynomial functions $\varphi_{r}(X)$ of $X \in \mathbf{C}_{n}^{n}$ by

$$
\begin{equation*}
\operatorname{det}\left(T 1_{n}-X\right)=\sum_{r=0}^{n}(-1)^{r} \varphi_{r}(X) T^{n-r} \tag{17.15}
\end{equation*}
$$

Notice that $\varphi_{\nu}(\operatorname{diag}[X, 0])=\varphi_{\nu}(X)$ for $\nu \leq n$. Now we have

$$
\begin{align*}
& \omega(2 \pi y, h ; \lambda+1,0)=2^{-p \lambda} \pi^{\iota p(n-p) / 2} \mathbf{e}^{n}(i h y) \delta_{+}(4 \pi h y)^{-1}  \tag{17.16}\\
& \cdot \sum_{\nu=0}^{p}(-1)^{\nu} b_{\nu}(\iota(n-p) / 2) \varphi_{p-\nu}(4 \pi h y) \text { if } q=0, \text { where } \\
& b_{0}(\alpha)=1 \text { and } b_{\nu}(\alpha)=\prod_{m=0}^{\nu-1}(\alpha+(\iota m / 2)) \text { if } \nu>0 .
\end{align*}
$$

This was proved in [S93, Lemma 9.2 (iv)]. Put $X=\pi h y$ in (17.15) and multiply by $(T \pi y)^{-1}$; then we obtain $\operatorname{det}\left[(\pi y)^{-1}-T^{-1} h\right]=\sum_{r=0}^{n} \operatorname{det}(\pi y)^{-1} \varphi_{r}(\pi h y)(-T)^{-r}$, which shows that
(17.17) $\operatorname{det}(\pi y)^{-1} \varphi_{r}(\pi h y)$ is a polynomial of degree $\leq n-r$ in the entries of $(\pi y)^{-1}$ with coefficients in the field generated over $\mathbf{Q}$ by the entries of $h$.
17.5. Lemma. Define $L_{\mathfrak{c}}(s, \psi)$ by (16.9) with a Hecke character $\psi$ of $F$.
(1) If $0 \geq m \in \mathbf{Z}, \psi_{v}\left(x_{v}\right)=\operatorname{sgn}\left(x_{v}\right)^{m}$ for some $v \in \mathbf{a}$, and $\mathfrak{c} \neq \mathfrak{g}$, then $L_{\mathfrak{c}}(m, \psi)=0$.
(2) If $\psi_{\mathbf{a}}(x)=\operatorname{sgn}\left(x_{\mathbf{a}}\right)^{m \mathbf{a}}$ with $0<m \in \mathbf{Z}$, then $L_{\mathbf{c}}(m, \psi) \in \pi^{d m} \mathbf{Q}_{\mathbf{a b}}$ and $L_{\mathrm{c}}(1-m, \psi) \in \mathbf{Q}_{\mathrm{ab}}$.
(3) If $\psi_{\mathbf{a}}(x)=\operatorname{sgn}\left(x_{\mathbf{a}}\right)^{m \mathbf{a}}$ with $0 \leq m \in \mathbf{Z}$, then $L_{\mathfrak{c}}(s, \psi)$ at $s=-m$ has a zero of order at least $d$ except when $m=0, \psi=1$, and $\mathfrak{c}=\mathfrak{g}$, in which case the order of zero is $d-1$.

Proof. Take $t \in \mathbf{Z}^{\mathbf{a}}$ so that $0 \leq t_{v} \leq 1$ for every $v \in \mathbf{a}$ and $m-t \in 2 \mathbf{Z}^{\mathbf{a}}$. Then $L(s, \psi) \prod_{v \in \mathbf{a}} \Gamma\left(\left(s+t_{v}\right) / 2\right)$ is holomorphic on $\mathbf{C}$ except for possible simple poles at $s=0$ and $s=1$, which occur if and only if $\psi=1$. Assertions (1) and (3) can easily be derived from this fact. Assertion (2) is included in a result which we prove as Theorem 18.12 in the next section.
17.6. Proposition. Let $E^{*}$ be defined as above with $\kappa=0, \mathfrak{c} \neq \mathfrak{g}$, and $k=$ $\mu \mathbf{a}, \mu \in 2^{-1} \mathbf{Z}$; suppose $\mu \geq \lambda$; let $h \in S$ and $r=\operatorname{rank}(h)$. In Case $S U$ let $\theta$ be the quadratic Hecke character of $F$ corresponding to $K / F$. Then $c_{h}(y, s)$ is finite at $s=\mu / 2$. Moeover, $c_{h}(y, \mu / 2) \neq 0$ only in the following cases, and in each case, except in case (viii), the value can be described as follows:
(I) $h=0$.
(i) $\quad \mu=\lambda: \quad c_{0}(y, \mu / 2) \in \mathbf{Q}_{\mathrm{ab}}$.
(ii) (Case SP) $\mu=(n+2) / 2, F=\mathbf{Q}$, and $\chi^{2}=1$ :

$$
c_{0}(y, \mu / 2)=c \operatorname{det}(y)^{-1 / 2} \text { with } c \in \mathbf{C} ; c \in \pi^{-\lambda} \mathbf{Q}_{\mathrm{ab}} \text { if } \mu \notin \mathbf{Z}
$$

(iii) (Case SP) $2<\mu=(n+3) / 2, F=\mathbf{Q}$, and $\chi^{2}=1$ :

$$
c_{0}(y, \mu / 2)=c \pi^{-n} \operatorname{det}(y)^{-1} \text { with } c \in \mathbf{Q}_{\mathrm{ab}}
$$

(iv) (Case SP) $\mu=2, n=1, F=\mathbf{Q}$, and $\chi=1$ :

$$
c_{0}(y, \mu / 2)=c \pi^{-1} y^{-1} \text { with } c \in \mathbf{Q}_{\mathrm{ab}} .
$$

(v) (Case SU) $\mu=n+1, F=\mathbf{Q}$, and $\chi=\theta^{n-1}$ :

$$
c_{0}(y, \mu / 2)=c \pi^{-n} \operatorname{det}(y)^{-1} \text { with } c \in \mathbf{Q}_{\mathrm{ab}} .
$$

(II) $h \neq 0$.
(vi) $\mu>\lambda$ and $h_{v}>0$ for every $v \in \mathbf{a}$ :

$$
c_{h}(y, \mu / 2)=c \mathbf{e}_{\mathbf{a}}^{n}(i h y) \text { with } c \in \mathbf{Q}_{\mathbf{a b}} .
$$

(vii) $\mu=\lambda$ and $h_{v} \geq 0$ for every $v \in \mathbf{a}$ :

$$
c_{h}(y, \mu / 2)=c \mathbf{e}_{\mathbf{a}}^{n}(i h y) \text { with } c \in \mathbf{Q}_{\mathrm{ab}} .
$$

(viii) (Case SP) $\mu=(n+2) / 2, F=\mathbf{Q}, \quad \chi^{2}=1$.
(ix) (Case SP) $\mu=(n+3) / 2, F=\mathbf{Q}, \chi^{2}=1, \quad 0<r<n$, and $h \geq 0$ :
$c_{h}(y, \mu / 2)=c \pi^{e} \operatorname{det}(\pi y)^{-1} \delta_{+}(\pi h y) \omega(2 \pi y, h ; \lambda+1,0)$ with $c \in \mathbf{Q}_{\mathrm{ab}}$ and $e=r(r-n) / 2$.
(x) (Case SU) $\mu=n+1, F=\mathbf{Q}, \chi=\theta^{n-1}, 0<r<n$, and $h \geq 0$ :
$c_{h}(y, \mu / 2)=c \pi^{e} \operatorname{det}(\pi y)^{-1} \delta_{+}(\pi h y) \omega(2 \pi y, h ; n+1,0)$ with $c \in \mathbf{Q}_{\mathrm{ab}}$ and $e=r(r-n)$.

Proof. By Propositions 16.9 and $16.10, c_{h}(y, s)$ in Case SP is easy finite nonvanishing factors times

$$
\begin{equation*}
f(s) \Lambda_{\mathfrak{c}}(s)^{-1} \Lambda_{h}(s) \prod_{v \in \mathbf{a}} \xi\left(y_{v}, h_{v} ; s+\mu / 2, s-\mu / 2\right) \tag{}
\end{equation*}
$$

where we denote by $f(s)$ the product $\prod_{v \in \mathbf{c}}$ of (16.46). Since we know the explicit forms of $\Lambda_{c}$ and $\Lambda_{h}$, we can derive our assertions from their properties and the formulas concerning $\xi$ in $\S 17.4$, combined with Lemma 17.5. The verification is fairly straightforward, but lengthy, as there are many combinations of $n, \mu, \chi$, and $h$, which produce different results. Therefore, we discuss here only a few typical cases in Case SP, leaving the remaining cases to the reader. Case SU can be handled in a similar and much simpler way by employing the formulas of [S97, Propositions 18.14 and 19.2].

First of all, since $\chi_{\mathbf{a}}(x)=\operatorname{sgn}\left(x_{\mathbf{a}}\right)^{\mu \mathbf{a}}$, from Lemma 17.5 (2) we see that

$$
\begin{align*}
0 & \neq \Lambda_{\mathrm{c}}(\mu / 2) \in \pi^{d f} \mathbf{Q}_{\mathrm{ab}}, \quad \text { where }  \tag{17.18a}\\
f & =\left\{\begin{array}{l}
\mu(2 m+1)-m(m+1) \text { with } m=[n / 2] \text { if } \mu \in \mathbf{Z}, \\
m(2 \mu-m) \text { with } m=[(n+1) / 2] \text { if } \mu \notin \mathbf{Z} .
\end{array}\right.
\end{align*}
$$

We also observe, for $(n+1) / 2 \leq \mu \in 2^{-1} \mathbf{Z}$, that

$$
\Gamma_{n}^{1}(\mu) \in \begin{cases}\pi^{\left[n^{2} / 4\right]} \mathbf{Q}^{\times} & \text {if } \quad \mu n \in \mathbf{Z}  \tag{17.18b}\\ \pi^{\left(n^{2}+1\right) / 4} \mathbf{Q}^{\times} & \text {if } \quad \mu n \notin \mathbf{Z}\end{cases}
$$

(I) Suppose $h=0$ and $\mu \in \mathbf{Z}$; then in ( ${ }^{*}$ ) we have $f=1$ and $\xi(y, 0 ; \cdots)$ is given by (17.13). Thus $\left(^{*}\right)$ takes the form

$$
\begin{align*}
& \Lambda_{\mathfrak{c}}(s)^{-1} L_{\mathfrak{c}}(2 s-n, \chi) \prod_{i=1}^{[n / 2]} L_{\mathfrak{c}}\left(4 s-2 n+2 i-1, \chi^{2}\right)  \tag{}\\
& \cdot a \cdot 2^{-d n s} \pi^{d n \lambda} \operatorname{det}(y)^{(\lambda-2 s) \mathbf{a}} \Gamma_{n}^{1}(2 s-\lambda)^{d} \Gamma_{n}^{1}(s+\mu / 2)^{-d} \Gamma_{n}^{1}(s-\mu / 2)^{-d}
\end{align*}
$$

with $a \in \mathbf{Q}_{\mathrm{ab}}$. Take $\mu=(n+2) / 2$ with even $n$, for example. By Lemma 17.5 (3), $L_{\mathrm{c}}(2 s-n, \chi)$ has a zero of order at least $d$ at $s=\mu / 2$ for any $\chi$. Now $\prod_{i=1}^{[n / 2]}$ is finite everywhere except when $\chi^{2}=1$, in which case the factor for $i=n / 2$ has a pole of order 1 at $s=\mu / 2$; the other factors belong to $\mathbf{Q}_{\mathrm{ab}}$. Now, from (17.18b) we easily see that the product of the last gamma factors of (**) at $s=\mu / 2$ produces a rational number times $\pi^{n^{2} / 4}$. Thus the whole product is 0 if $\chi^{2} \neq 1$ or $F \neq \mathbf{Q}$. If $\chi^{2}=1$ and $F=\mathbf{Q}$, we obtain a constant times $\operatorname{det}(y)^{-1 / 2}$, as stated in (ii).

Still with $\mu=(n+2) / 2$, assume that $n \notin 2 \mathbf{Z}$. Then the first line of $\left({ }^{* *}\right)$ has a different expression, but the second line is the same. This time $\Gamma_{n}^{1}(2 s-\lambda)^{d} \Gamma_{n}^{1}(s-$ $\mu / 2)^{-d}$ has a zero of order $d$ at $s=\mu / 2$. Then we obtain (ii) in this case in a similar manner.

Next suppose $\mu=(n+3) / 2$ with odd $n$. Again we see that the nonvanishing can occur only if $F=\mathbf{Q}$ and $\chi^{2}=1$, in which case

$$
\Lambda_{0}(s)=L_{\mathfrak{c}}(2 s-n, \chi) \prod_{i=1}^{[n-1) / 2} L_{\mathfrak{c}}(4 s-2 n+2 i-1,1)
$$

Suppose $n=1$; then $\mu=2$ and we easily see that the nonvanishing can occur only when $\chi=1$, and the result is as stated in (iv). If $n>1$, we have $\mu<n$, and so $L_{\mathrm{c}}(\mu-n, \chi) \in \mathbf{Q}_{\mathrm{ab}}$; also $\prod_{i=1}^{(n-1) / 2}$ has a pole of order 1 at $s=\mu / 2$ and its residue is a rational number times $\prod_{i=1}^{(n-3) / 2} L_{\mathfrak{c}}(1-2 i, 1)$, which belongs to $\mathbf{Q}_{\mathrm{ab}}$. In this way we obtain (iii) and (iv).
(II) Suppose $h \neq 0$; define $\rho_{h}$ as in Proposition 16.10 and put $r=\operatorname{rank}(h)$; let $p_{v}$ (resp. $q_{v}$ ) be the number of positive (resp. negative) characteristic roots of $h_{v}$. Then $\left(\chi \rho_{h}\right)_{v}(x)=\operatorname{sgn}\left(x_{v}\right)^{[\mu]+[r / 2]+q_{v}}$ for $v \in \mathbf{a}$. First suppose $r=n$ and $2 \mu+n \notin 2 \mathbf{Z}$; then $\Lambda_{h}=1$. From (17.11) we see that the last factor $\prod_{v \in \mathbf{a}}$ of $\left({ }^{*}\right)$ is finite at $s=\mu / 2$, and is nonzero only if $q_{v}=0$ for every $v$, in which case (17.12) shows that $c_{h}(y, \mu / 2)$ is an element of $\mathbf{Q}_{\mathrm{ab}}$ times $\Lambda_{\mathfrak{c}}(\mu / 2)^{-1} \Gamma_{n}^{1}(\mu)^{-d} \pi^{d n \mu} \mathbf{e}_{\mathbf{a}}^{n}(i h y)$. Employing (17.18a, b), we see that $c_{h}(y, \mu / 2)=c \mathbf{e}_{\mathbf{a}}^{n}(i h y)$ with $c \in \mathbf{Q}_{\mathrm{ab}}$.

The case $0<r<n$ is more complicated. Suppose $\mu \notin \mathbf{Z}, r<n \in 2 \mathbf{Z}$, and $r \notin 2 \mathbf{Z}$, for example; then
$\left({ }^{* * *}\right) \quad \Lambda_{h}(s)=L_{\mathrm{c}}\left(2 s-n+r / 2, \chi \rho_{h}\right) \prod_{i=1}^{[(n-r) / 2]} L_{\mathrm{c}}\left(4 s-2 n+r+2 i-1, \chi^{2}\right)$.
By (17.11), $\prod_{v \in \mathbf{a}} \xi(\cdots)$ of (*) is a finite factor times $\prod_{v \in \mathbf{a}} \Gamma_{t}^{1}(2 \sigma+\mu-\lambda) \Gamma_{t+q_{v}}^{1}(\sigma)^{-1}$ with $\sigma=s-\mu / 2$, where $t=n-r$. First suppose $\mu=\lambda$; then $t$ is odd and ( ${ }^{* * *)}$ is finite at $s=\mu / 2$; the last $\prod_{v \in \mathbf{a}}$ is 0 at $\sigma=0$ if $q_{v} \geq 2$ for some $v \in \mathbf{a}$. Thus we may assume that $q_{v} \leq 1$ for every $v \in \mathbf{a}$. Suppose $q_{v}=1$ for some $v$. Then
$[\mu]+[r / 2]+q_{v}=(n+r+1) / 2 \equiv \mu-n+r / 2(\bmod 2)$. Therefore the first factor of $\left({ }^{* * *}\right)$ is 0 at $s=\mu / 2$ by Lemma 17.5 (1). Thus the nonvanishing occurs only when $q_{v}=0$ for every $v \in \mathbf{a}$, that is, when $h$ is totally nonnegative. Then employing (17.14), we obtain the desired formula as given in (vii).

Next suppose $\mu=\lambda+1$. Then ( ${ }^{* * *}$ ) is finite if $\chi^{2} \neq 1$; it has a pole of order at most 1 at $s=\mu / 2$ if $\chi^{2}=1$. This time the product of gamma functions has a zero of order at least $d$ at $\sigma=0$. Therefore the nonvanishing can occur only if $\chi^{2}=1$, $F=\mathbf{Q}$, and $\Lambda_{h}$ has a pole of order 1 at $s=\mu / 2$. In that special case suppose $q>1$; then $\Gamma_{t}^{1}(2 \sigma+1) / \Gamma_{t+q}^{1}(\sigma)$ has a zero of order at least 2 at $\sigma=0$, so that the nonvanishing can occur only if $q \leq 1$. If $r=n-1$, then $\chi=\rho_{h}$ and so the signature formula for $\chi \rho_{h}$ says that $q \in 2 \mathbf{Z}$, a contradiction. Thus we may assume that $q \leq 1$ and $r \leq n-3$. Suppose $q=1$. Then $[\mu]+[r / 2]+q=(n+3+r) / 2 \equiv \mu-n+r / 2$ $(\bmod 2)$. Therefore $L_{\mathfrak{c}}\left(2 s-n+r / 2, \chi \rho_{h}\right)$ is 0 at $s=\mu / 2$ by Lemma 17.5 (1). Thus we may assume $q=0$. Then from (17.11) we obtain the desired formula as given in (ix). All the remaining cases can be handled more or less in the same manner.

To make our statements shorter, we make the following convention: whenever we speak of a function $f(z, \mu / 2)$ belonging to a set, it means that $f(z, s)$ is finite at $s=\mu / 2$, and the value as a function of $z$ belongs to the set in question; whenever we speak of $E(z, s ; \nu \mathbf{a}, \ldots)$ or $D(z, s ; \nu \mathbf{a}, \ldots)$, we assume that $\nu \in 2^{-1} \mathbf{Z}$ in Case SP and $\nu \in \mathbf{Z}$ in case SU . In $\S 14.11$ we defined $\mathcal{N}_{\omega}^{p}(W)$ for $p \in \mathbf{Z}^{\mathbf{a}}$, a representation $\omega$ of $G L_{n}(\mathbf{C})^{\mathbf{b}}$, and a subfield $W$ of $\mathbf{C}$. In Case SU, viewing $\mathbf{Z}^{\mathbf{a}}$ as a submodule of $\mathbf{Z}^{\mathbf{b}}$ in an obvious way, we can speak of $\mathcal{N}_{k}^{p}(W)$ for $k \in \mathbf{Z}^{\mathbf{a}}$.
17.7. Theorem (Cases SP and SU). (i) If $\mu \geq \lambda$, then $E(z, \mu / 2 ; \mu \mathbf{a}, \Gamma)$ belongs to $\mathcal{M}_{\mu \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ except when $F=\mathbf{Q}$ and $\lambda+(1 / 2) \leq \mu \leq \lambda+1$.
(ii) If $F=\mathbf{Q}$ and $\mu=\lambda+1$, then $E(z, \mu / 2 ; \mu \mathbf{a}, \Gamma)$ belongs to $\mathcal{N}_{\mu}^{n}\left(\mathbf{Q}_{\mathrm{ab}}\right)$.
(iii) If $\mu \geq \lambda$, then $E(z, \mu / 2 ; \mu \mathbf{a}, \chi, \mathfrak{c})$ belongs to $\mathcal{M}_{\mu \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ except in the following four cases:
(A) Case SP: $\mu=(n+2) / 2, F=\mathbf{Q}$, and $\chi^{2}=1$;
(B) Case SP: $n=1, \mu=2, F=\mathbf{Q}$, and $\chi=1$;
(C) Case SP: $n>1, \mu=(n+3) / 2, F=\mathbf{Q}$, and $\chi^{2}=1$;
(D) Case SU: $\mu=n+1, F=\mathbf{Q}$, and $\chi=\theta^{n+1}$.
(iv) In Cases (B), (C), (D), $E(z, \mu / 2 ; \mu \mathbf{a}, \chi, \mathfrak{c})$ belongs to $\mathcal{N}_{\mu}^{n}\left(\mathbf{Q}_{\mathrm{ab}}\right)$.
(v) Taking an integer $\mu \leq \lambda$, put $\nu=2 \lambda-\mu$ and

$$
e= \begin{cases}n(n+2) / 4 & (\text { Case SP, } n \in 2 \mathbf{Z}) \\ (n+1) \lambda / 2 & \text { (all othe cases) }\end{cases}
$$

Then $D(z, \mu / 2 ; \nu \mathbf{a}, \chi, \mathfrak{c})$ belongs to $\pi^{d e} \mathcal{M}_{\nu \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ except in the following three cases:
(E) Case SP: $\mu=0, \mathfrak{c}=\mathfrak{g}$, and $\chi=1$;
(F) Case SP: $0<\mu \leq n / 2, \mathfrak{c}=\mathfrak{g}$, and $\chi^{2}=1$;
(G) Case SU: $0 \leq \mu<n, \mathfrak{c}=\mathfrak{g}$, and $\chi=\theta^{\mu}$.
(vi) If $n=1$ and $\mathfrak{c}=\mathfrak{g}$, then $D(z, 0 ; 2 \mathbf{a}, \chi, \mathfrak{c})$ belongs to $\pi^{d} \mathcal{M}_{2 \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ except when $F=\mathbf{Q}$, in which case it belongs to $\pi \mathcal{N}_{2}^{1}\left(\mathbf{Q}_{\mathrm{ab}}\right)$.

Proof. We first prove (iii) and (iv) when $\mathfrak{c} \neq \mathfrak{g}$. By Theorem 7.11, $\mathcal{M}_{\mu \mathrm{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ is stable under $f \mapsto f \| \alpha$ for every $\alpha \in G_{1}$ or $\alpha \in \mathcal{G}$. As for $\mathcal{N}_{\mu \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$, we need it only when $F=\mathbf{Q}$. Therefore, by Theorem 14.13, it is stable under $f \mapsto f \| \alpha$ for every such $\alpha$ (see $\S 14.14$ ). Thus it is sufficient to prove (iii) and (iv) for $E^{*}(z, s)$ instead
of $E(z, s ; \cdots)$. Now the desired facts for $E^{*}$ follow immediately from Proposition 17.6. The case in which $\mu>\lambda+1$ is the easiest: $E^{*}(z, \mu / 2)=\sum_{h \in S} b_{h} \mathbf{e}_{\mathbf{a}}^{n}(h z)$ with $b_{h} \neq 0$ only if $h$ is totally positive and $b_{h} \in \mathbf{Q}_{\mathrm{ab}}$. All other cases are similar. Special care must be taken when $\mu=\lambda+1$ and $F=\mathbf{Q}$. Here let us discuss only the case in which $n>1$ in Case SP. We assume $\chi^{2}=1$, since the function is holomorphic otherwise. By (iii) and (vi) of Proposition 17.6, $c_{0}(y, \mu / 2)=a \operatorname{det}(\pi y)^{-1}$ with $a \in \mathbf{Q}_{\mathrm{ab}}$ and $c_{h}(y, \mu / 2)=b \mathbf{e}_{\mathbf{a}}^{n}(i h y)$ with $b \in \mathbf{Q}_{\mathrm{ab}}$ if $\operatorname{det}(h) \neq 0$. For $0<\operatorname{rank}(h)=$ $r<n, c_{h}(y, \mu / 2)$, if nonzero, can be given by the formula of (ix) of Proposition 17.6. From (17.16) and (17.17) we easily see that $c_{h}(y, \mu / 2)=q\left((\pi y)^{-1}\right) \mathbf{e}_{\mathbf{a}}^{n}(i h y)$ with a $\mathbf{Q}_{\mathrm{ab}}$-rational polynomial $q$ of degree $\leq n-r$. Thus $E^{*}(z, \mu / 2)$ in this case belongs to $\mathcal{N}_{\mu \mathbf{a}}^{n}\left(\mathbf{Q}_{\mathrm{ab}}\right)$.

Next suppose $\mathfrak{c}=\mathfrak{g}$; then $\mu \in \mathbf{Z}$. Fix an arbitrary prime ideal $\mathfrak{p}$ of $F$. By Lemma 17.2 (3), $E(z, s ; k, \chi, \mathfrak{g})=\sum_{\alpha \in A} E(z, s ; k, \chi, \mathfrak{p}) \|_{k} \alpha$ with a finite subset $A$ of $\Gamma_{0}(\mathfrak{c})$. Therefore we obtain (iii) and (iv) in the case $\mathfrak{c}=\mathfrak{g}$ from those in the case $\mathfrak{c} \neq \mathfrak{g}$.

Once (iii) and (iv) are established, (i) and (ii) can be obtained by combining (iii) and (iv) with Lemma 17.2 (2), because of the stability of $\mathcal{M}_{\mu \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ and $\mathcal{N}_{\mu}^{n}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ under $\|_{\mu \mathbf{a}} \alpha$ as mentioned at the beginning.

As for (v), the method of proof is the same. In Case $\mathrm{SU}, \Lambda_{\mathrm{c}}$ is given in [S97, Proposition 19.2], which is the same as the product of $L$-functions in (17.7). If $\mathfrak{c} \neq \mathfrak{g}$, we again reduce the problem to $\Lambda_{\mathfrak{c}}(s) E^{*}(z, s ; \nu \mathbf{a})$. Thus the question is the behavior of $\Lambda_{\mathbf{c}}(s) c_{h}(y, s)$ at $s=\mu / 2$ with $k=\nu \mathbf{a}$. The analysis of the value is similar to that of $c_{h}(y, s)$ at $s=\mu / 2$ in the proof of Proposition 17.6, and so we do not go into details here. We note only that we need tha value $\xi(y, h ; \lambda, \mu-\lambda)$, which can be obtained from [S82, (4.34.K), (4.35.K)]. If $\mathfrak{c}=\mathfrak{g}$, with any fixed prime ideal $\mathfrak{p}$ we have

$$
\begin{gathered}
\qquad B(s) D(z, s ; k, \chi, \mathfrak{g})=\sum_{\alpha \in A} D(z, s ; k, \chi, \mathfrak{p}) \|_{k} \alpha \\
\text { with } \quad B(s)= \begin{cases}\left(1-\chi^{*}(\mathfrak{p}) N(\mathfrak{p})^{-2 s}\right) \prod_{i=1}^{[n / 2]}\left(1-\chi^{*}(\mathfrak{p})^{2} N(\mathfrak{p})^{2 i-4 s}\right) & \text { (Case SP), } \\
\prod_{i=0}^{n-1}\left(1-\left(\chi \theta^{i}\right)^{*}(\mathfrak{p}) N(\mathfrak{p})^{i-2 s}\right) & \text { (Case SU). }\end{cases}
\end{gathered}
$$

Observe that we can choose $\mathfrak{p}$ so that $B(\mu / 2) \neq 0$ if we exclude Cases (E), (F), and (G). Therefore we can derive the desired conclusion of (v) for $\mathfrak{c}=\mathfrak{g}$ from that for $\mathfrak{c}=\mathfrak{p}$.

We shall prove (vi) in §A2.14.
17.8. Theorem (Cases SP and SU). Suppose $0<\mu<\lambda$; put $s_{\mu}=\lambda-(\mu / 2)$. Then $E(z, s ; \mu \mathbf{a}, \chi, \mathfrak{c})$ has at most a simple pole at $s=s_{\mu}$, which occurs only if $\chi^{2}=1$ in Case $S P$ and $\chi=\theta^{\mu}$ in Case SU. The residue at $s=s_{\mu}$ is of the form $\pi^{-d \gamma} A \cdot g(z)$ with an element $g$ of $\mathcal{M}_{\mu \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ and constants $\gamma$ and $A$ given by

$$
\gamma= \begin{cases}n(n / 2-\mu) & (\text { Case SP, } \mu \in \mathbf{Z}, n \in 2 \mathbf{Z}) \\ (n-1)(n / 2-\mu)-1 & (\text { Case SP, } \mu \in \mathbf{Z}, n \notin 2 \mathbf{Z}) \\ (n+1-2 \mu)[(n+1) / 2)] & (\text { Case SP, } \mu \notin \mathbf{Z}), \\ n(n-\mu) & (\text { Case SU) }\end{cases}
$$

$$
A=R_{F} \cdot \begin{cases}L(n+1-\mu, \chi) \prod_{i=1}^{[n / 2-\mu]} \zeta_{F}(2 i+1) & (\text { Case SP, } \mu \in \mathbf{Z}) \\ \prod_{i=1}^{[n / 2-\mu]} \zeta_{F}(2 i+1) & \text { (Case SP, } \mu \notin \mathbf{Z}) \\ \prod_{i=2}^{n-\mu} L\left(i, \theta^{i-1}\right) & \text { (Case SU), }\end{cases}
$$

where $R_{F}$ is the regulator of $F$, and $\zeta_{F}$ is the Dedekind zeta function of $F$. Moreover, $g(z)=\sum_{h} a(h) \mathbf{e}_{\mathbf{a}}^{n}(h z)$ with $a(h)=0$ whenever $\operatorname{rank}(h)>2 \mu$ in Case SP and $\operatorname{rank}(h)>\mu$ in Case $S U$.

Proof. The method of proof is the same as in Theorem 17.7. Namely, we reduce the problem to $E^{*}(z, s)$, and study the behavior of the Fourier coefficients at $s=s_{\mu}$, employing the explicit form of $c(h, q, s)$ given in Proposition 16.9. In particular, we have to analyze $\xi(g, h ; s+\mu / 2, s-\mu / 2)$ at $s=s_{\mu}$, for which we need (17.11) and [S82, (4.35.K)]. The details may be left to the reader. There is one more nontrivial point: that the property $a(h)=0$ for $\operatorname{rank}(h)>t$ with a fixed $t$ can be preserved by the transformation $g \mapsto g \| \alpha$ for every $\alpha \in G_{1}$. For this, see [S94b, (5.14)].
17.9. Theorem (Cases SP and SU). Let $\Phi$ be the Galois closure of $K$ over $\mathbf{Q}$ and let $k$ be a weight (which means that $k \in \mathbf{Z}^{\mathbf{a}}$ in Case SU); suppose that $k_{v} \geq \lambda$ for every $v \in \mathbf{a}$ and $k_{v}-k_{v^{\prime}} \in 2 \mathbf{Z}$ for every $v, v^{\prime} \in \mathbf{a}$.
(ia) Let $\mu$ be an element of $2^{-1} \mathbf{Z}$ such that $\lambda \leq \mu \leq k_{v}$ and $\mu-k_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$. Then $E(z, \mu / 2 ; k, \Gamma)$ and $E(z, \mu / 2 ; k, \chi, \mathfrak{c})$ belong to $\pi^{\alpha} \mathcal{N}_{k}^{t}\left(\Phi \mathbf{Q}_{\mathrm{ab}}\right)$ except when $F=\mathbf{Q}$ and $\mu=(n+2) / 2$ in Case $S P$, where $\alpha=(n / 2) \sum_{v \in \mathbf{a}}\left(k_{v}-\mu\right)$ and

$$
t= \begin{cases}n(k-\mu+2) / 2 & \text { if } \mu=\lambda+1 \quad \text { and } \quad F=\mathbf{Q} \\ n(k-\mu \mathbf{a}) / 2 & \text { otherwise }\end{cases}
$$

(ib) For $\mu$ and $\alpha$ as above, $E(z, \mu / 2 ; k, \chi, \mathrm{c})$ belongs to $\pi^{\alpha} \mathcal{N}_{k}^{u}\left(\Phi \mathbf{Q}_{\mathrm{ab}}\right)$ except in Case (A) of Theorem 17.7, where

$$
u= \begin{cases}n(k-\mu+2) / 2 & \text { in Cases }(B),(C),(D) \text { of Theorem 17.7 } \\ n(k-\mu \mathbf{a}) / 2 & \text { otherwise }\end{cases}
$$

(ii) Let $\mu \in 2^{-1} \mathbf{Z}$; suppose that $2 \lambda-k_{v} \leq \mu \leq k_{v}$ and $|\mu-\lambda|+\lambda-k_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$. Suppose also that $\mu, F, \chi$, and $\mathfrak{c}$ do not fall into Cases (A), (E), $(F)$, and (G) of Theorem 17.7. Then $D(z, \mu / 2 ; k, \chi, \mathfrak{c})$ belongs to $\pi^{\beta} \mathcal{N}_{k}^{r}\left(\Phi \mathbf{Q}_{\mathrm{ab}}\right)$, where

$$
r= \begin{cases}n(k-\mu+2) / 2 & \text { in Cases }(B),(C),(D) \\ (n / 2)(k-|\mu-\lambda| \mathbf{a}-\lambda \mathbf{a}) & \text { otherwise }\end{cases}
$$

and $\beta=(n / 2) \sum_{v \in \mathbf{a}}\left(k_{v}+\mu\right)-d e$ with

$$
e= \begin{cases}{\left[(n+1)^{2} / 4\right]-\mu} & (\text { Case SP: } 2 \mu+n \in 2 \mathbf{Z} \text { and } \mu \geq \lambda), \\ {\left[n^{2} / 4\right]} & (\text { Case SP: } 2 \mu+n \notin 2 \mathbf{Z} \text { or } \mu<\lambda), \\ n(n-1) / 2 & (\text { Case SU). }\end{cases}
$$

(iii) Suppose $n=1, \mu=0$, and $2 \leq k_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$. Then $D(z, 0 ; k, \chi, \mathbf{c})$ belongs to $\pi^{\beta} \mathcal{N}_{k}^{r}\left(\Phi \mathbf{Q}_{\mathbf{a b}}\right)$ with $r=(k-2 \mathbf{a}) / 2$ and $\beta=\sum_{v \in \mathbf{a}} k_{v} / 2$,
except when $F=\mathbf{Q}$ and $\mathfrak{c}=\mathbf{Z}$. If $F=\mathbf{Q}$ and $\mathfrak{c}=\mathbf{Z}$, then it belongs to $\pi^{k / 2} \mathcal{N}_{k}^{k / 2}\left(\mathbf{Q}_{\mathrm{ab}}\right)$.

Assertion (iii) means that if $n=1$ we have results even in Cases (E) and (G). It is conjecturable that similar results hold in Cases (E), (F), and (G) even for $n>1$.

Proof. For $p \in \mathbf{Z}^{\mathbf{a}}$ and a weight $q$ define the operator $\Delta_{q}^{p}$ by $\Delta_{q}^{p} f=\left(D_{\rho}^{Z} f\right)(\psi)$ with $\rho(a, b)=\operatorname{det}(b)^{q}, Z=\bigotimes_{v \in \mathbf{a}} Z_{v}=\mathbf{C} \psi \subset S_{n p}(T)$, where $Z_{v}=\mathbf{C} \psi_{v} \subset S_{n p_{v}}\left(T_{v}\right)$, $\psi=\prod_{v \in \mathbf{a}} \psi_{v}$, and $\psi_{v}=\operatorname{det}(x)^{p_{v}}$. Then, by (12.24c) and Theorem 14.12 (4),

$$
\begin{equation*}
\Delta_{q}^{p} \mathcal{N}_{q}^{t}\left(\Phi \mathbf{Q}_{\mathrm{ab}}\right) \subset \pi^{n|p|} \mathcal{N}_{q+2 p}^{t+n p}\left(\Phi \mathbf{Q}_{\mathrm{ab}}\right) \tag{17.19}
\end{equation*}
$$

(See $\S 14.14$ if $q \notin \mathbf{Z}^{\mathbf{a}}$.) Let simply $E(z, s ; k)$ denote $E(z, s ; k, \Gamma)$ or $E(z, s ; k, \chi, \mathfrak{c})$. We now apply Lemma 13.9 to each term of $E$ with $\zeta(x)=\operatorname{det}(x)^{p_{v}}$. Employing the formula for $\psi_{Z}$ given in Theorem 12.13, we find that

$$
\begin{equation*}
\Delta_{q}^{p} E(z, s ; q)=c_{q}^{p}(s)(i / 2)^{n|p|} E(z, s ; q+2 p) \tag{17.20}
\end{equation*}
$$

$$
\text { with } c_{q}^{p}(s)= \begin{cases}\prod_{v \in \mathbf{a}} \prod_{a=1}^{n} \prod_{b=1}^{p_{v}}\left\{-s-\left(q_{v} / 2\right)-b+(a+1) / 2\right\} & \text { (Case SP), } \\ \prod_{v \in \mathbf{a}} \prod_{a=1}^{n} \prod_{b=1}^{p_{v}}\left\{-s-\left(q_{v} / 2\right)-b+a\right\} & \text { (Case SU). }\end{cases}
$$

This is so even when $q \notin \mathbf{Z}^{\mathbf{a}} ;$ see $\S 14.14$. Now, given $\mu$ as in (ia), put $p=(k-\mu \mathbf{a}) / 2$. Then $0 \leq p \in \mathbf{Z}^{\mathbf{a}}$, and (17.20) with $q=\mu \mathbf{a}$ yields

$$
\begin{equation*}
\Delta_{\mu \mathbf{a}}^{p} E(z, \mu / 2 ; \mu \mathbf{a})=c_{\mu \mathbf{a}}^{p}(\mu / 2)(i / 2)^{n|p|} E(z, \mu / 2 ; k) . \tag{17.21}
\end{equation*}
$$

Observing that $c_{\mu \mathbf{a}}^{p}(\mu / 2) \neq 0$, we obtain our assertions of (ia) and (ib) from (i) and (ii) of Theorem 17.7 and (17.19).

Next, let $\mu$ be given as in (ii). If $\mu \geq \lambda$, the above proof is applicable also to this case, and we obtain the desired result from (iii) and (iv) of Theorem 17.7, except that we have to multiply the functions by the product of certain values of $L$ functions, and hence a power of $\pi$ appears as stated. If $\mu<\lambda$, we put $\nu=2 \lambda-\mu$ and $p=(k-\nu \mathbf{a}) / 2$. Then $\nu>\lambda, 0 \leq p \in \mathbf{Z}^{\mathbf{a}}$, and

$$
\begin{equation*}
\Delta_{\nu \mathbf{a}}^{p} D(z, \mu / 2 ; \nu \mathbf{a}, \chi, \mathfrak{c})=(i / 2)^{n|p|} c_{\nu \mathbf{a}}^{p}(\mu / 2) D(z, \mu / 2 ; k, \chi, \mathfrak{c}) \tag{17.22}
\end{equation*}
$$

Observing that $c_{\nu \mathbf{a}}^{p}(\mu / 2) \neq 0$, we obtain our assertion of (ii) from (v) of Theorem 17.7 and (17.19). Assertion (iii) follows from (vi) of Theorem 17.7 in the same manner.
17.10. In our later investigations we need the series $E_{\mathbf{A}}$ of (16.27) in Case UT. We put $E^{*}(x, s)=E_{\mathbf{A}}\left(x \eta_{\mathrm{h}}^{-1}, s\right)$. Given $r \in G_{\mathrm{h}}$ in Case UT, we define functions $E_{r}(z, s)$ and $E_{r}^{*}(z, s)$ of $(z, s) \in \mathcal{H} \times \mathbf{C}$ by

$$
\begin{align*}
& E_{r}(x(\mathbf{i}), s)=E_{r}(x(\mathbf{i}), s ; k, \chi, \mathfrak{c})=j_{x}^{k}(\mathbf{i}) E_{\mathbf{A}}(r x, s)  \tag{17.23a}\\
& E_{r}^{*}(x(\mathbf{i}), s)=E_{r}^{*}(x(\mathbf{i}), s ; k, \chi, \mathfrak{c})=j_{x}^{k}(\mathbf{i}) E_{\mathbf{A}}^{*}(r x, s)
\end{align*}
$$

for $x \in G_{\mathbf{a}}$. Then we put

$$
\begin{equation*}
D_{r}(z, s ; k, \chi, \mathfrak{c})=E_{r}(z, s ; k, \chi, \mathfrak{c}) \prod_{i=0}^{n-1} L_{\mathfrak{c}}\left(2 s-i, \chi_{1} \theta^{i}\right) \tag{17.24}
\end{equation*}
$$

where $\chi_{1}$ is the restriction of $\chi$ to $F_{\mathbf{A}}^{\times}$. Here recall that in Case UT, $\chi$ is a Hecke character of $K$.
17.11. Lemma. Let $\chi$ be a Hecke character of $K$ satisfying (16.24a) with $\ell \in$ $\mathbf{Z}^{\mathbf{a}}$ and $\kappa=0$. Then $\chi(c)$ for every $c \in K_{\mathbf{h}}^{\times}$and $\chi^{*}(\mathfrak{a})$ for every $\mathfrak{r}$-ideal $\mathfrak{a}$ are algebraic.

Proof. Since $\chi^{*}(\mathfrak{a})$ is 0 or $\chi(c)$ for some $c \in K_{\mathbf{h}}^{\times}$, it is sufficient to treat $\chi(c)$. Given $c$, we can find $\alpha \in K^{\times}$and a positive integer $m$ such that $c^{m} \mathfrak{r}=\alpha \mathfrak{r}$. Then $c^{m} / \alpha=e f$ with $e \in \prod_{v \in \mathbf{h}} \mathbf{r}_{v}^{\times}$and $f \in K_{\mathbf{a}}^{\times}$. We can find an integral $\mathfrak{g}$ ideal $\mathfrak{c}$ with which (16.24b) is satisfied. We can also find a positive integer $\nu$ such that $e_{v}^{\nu}-1 \in \mathfrak{r}_{v} \boldsymbol{c}_{v}$ for every $v \mid \mathfrak{c}$. Then $\chi\left(e^{\nu}\right)=1$, and hence $\chi(e)$ is a root of unity. Now $f=\alpha_{\mathbf{a}}^{-1}$, so that $\chi(f)=|\alpha|^{\ell} \alpha^{-\ell} \in \overline{\mathbf{Q}}$. Since $\chi(\alpha)=1$, we have $\chi(c)^{m}=\chi(e) \chi(f) \in \overline{\mathbf{Q}}$, which proves our lemma.

It should be noted that given $\ell \in \mathbf{Z}^{\mathbf{a}}$, a Hecke character $\chi$ as in the above lemma always exists; see [S97, Lemma 11.14 (3)].
17.12. Theorem (Case UT). Let $k \in \mathbf{Z}^{\mathbf{b}}$ and $\mu \in \mathbf{Z}$; put $m=\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$. Let $\chi$ be a Hecke character of $K$ satisfying (16.24a, b) with $\kappa=0$, and $\chi_{1}$ the restriction of $\chi$ to $F_{\mathbf{A}}^{\times}$. Further let $K^{\prime}$ be the reflex field (in the sense of §9.4) of the CM-type $\left(K,\left\{\tau_{v}\right\}\right)$ we fixed in $\S 3.5$, and $K_{\chi}$ the field generated over $K^{\prime}$ by the values $\chi(c)$ for all $c \in K_{\mathbf{h}}^{\times}$. (Then $K_{\chi} \subset \overline{\mathbf{Q}}$ by Lemma 17.11.) Then the following assertions hold:
(i) Suppose $m_{v}=\mu \geq n$ for every $v \in \mathbf{a}$; then $E_{r}(z, \mu / 2 ; k, \chi, \mathfrak{c})$ belongs to $\mathcal{M}_{k}\left(K_{\chi} \mathbf{Q}_{\mathrm{ab}}\right)$ except when $\mu=n+1, F=\mathbf{Q}$, and $\chi_{1}=\theta^{n+1}$.
(ii) Let $k$ and $\mu$ be as in (i). If $\mu=n+1, F=\mathbf{Q}$, and $\chi_{1}=\theta^{n+1}$, then $E_{r}(z, \mu / 2 ; k, \chi, \mathfrak{c})$ belongs to $\mathcal{N}_{k}^{n}\left(K_{\chi} \mathbf{Q}_{\mathrm{ab}}\right)$.
(iii) Let $k$ and $\mu$ be as in (i); put $\nu=2 n-\mu$ and $e=n(n+1) / 2$. Then $D_{r}(z, \nu / 2 ; k, \chi, \mathfrak{c})$ belongs to $\pi^{d e} \mathcal{M}_{k}\left(K_{\chi} \mathbf{Q}_{\mathrm{ab}}\right)$ except when $0 \leq \nu<n, \mathfrak{c}=\mathfrak{g}$, and $\chi_{1}=\theta^{\mu}$.
(iv) Let $\Phi$ be the Galois closure of $K$ over $\mathbf{Q}$. Suppose $n \leq \mu \leq m_{v}$ and $\mu-m_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$. Then $E_{r}(z, \mu / 2 ; k, \chi, \mathfrak{c})$ belongs to $\pi^{\alpha} \mathcal{N}_{k}^{t}\left(\Phi K_{\chi} \mathbf{Q}_{\mathrm{ab}}\right)$, where $\alpha=(n / 2) \sum_{v \in \mathbf{a}}\left(m_{v}-\mu\right)$ and

$$
t= \begin{cases}n(m-\mu+2) / 2 & \text { if } \mu=n+1, F=\mathbf{Q}, \quad \text { and } \quad \chi_{1}=\theta^{n+1} \\ n(m-\mu \mathbf{a}) / 2 & \text { otherwise. }\end{cases}
$$

(v) Suppose $2 n-m_{v} \leq \mu \leq m_{v}$ and $m_{v}-\mu \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$. Then $D_{r}(z, \mu / 2 ; k, \chi, \mathfrak{c})$ belongs to $\pi^{\beta} \mathcal{N}_{k}^{t}\left(\Phi K_{\chi} \mathbf{Q}_{\mathrm{ab}}\right)$, except when $0 \leq \mu<n, \mathfrak{c}=\mathfrak{g}$, and $\chi_{1}=\theta^{\mu}$, where $\beta=(n / 2) \sum_{v \in \mathbf{a}}\left(m_{v}+\mu\right)-d n(n-1) / 2$ and

$$
t=\left\{\begin{array}{l}
n(m-\mu+2) / 2 \text { if } \mu=n+1, F=\mathbf{Q}, \quad \text { and } \quad \chi_{1}=\theta^{n+1}, \\
(n / 2)(m-|\mu-n| \mathbf{a}-n \mathbf{a}) \quad \text { otherwise. }
\end{array}\right.
$$

REMARK. (1) If $k_{v \rho}=k_{v}$ for every $v \in \mathbf{a}$, then $\chi_{\mathbf{a}}=1$, so that $\chi$ is of finite order. Therefore $K_{\chi} \subset \mathbf{Q}_{\mathrm{ab}}$ in that case.
(2) We can define $E(z, s ; k, \Gamma)$ for $k \in \mathbf{Z}^{\mathbf{a}}$ and a congruence subgroup $\Gamma$ of $U\left(\eta_{n}\right)$ in Case UT by (17.2), (17.3), and (17.3a). Then (17.5) holds for $\Gamma_{1} \subset \Gamma \subset$ $U\left(\eta_{n}\right)$. Now we can take $\Gamma_{1}$ in $S U\left(\eta_{n}\right)$. Therefore the nature of $E(z, \mu / 2 ; k, \Gamma)$ for $\Gamma \subset U\left(\eta_{n}\right)$ can be reduced to the case $\Gamma \subset S U\left(\eta_{n}\right)$, which is given in Theorems 17.7 and 17.9.

Proof. We first note that $E_{\mathbf{A}}^{*}$ has a Fourier expansion of the form

$$
E_{\mathbf{A}}^{*}\left(\left[\begin{array}{cc}
q & \sigma \widehat{q}  \tag{17.25}\\
0 & \widehat{q}
\end{array}\right]\right)=\sum_{h \in S} c(h, q, s) \mathbf{e}_{\mathbf{A}}^{n}(h \sigma) \quad\left(q \in G L_{n}(K)_{\mathbf{A}}, \sigma \in S_{\mathbf{A}}\right)
$$

with $c(h, q, s) \in \mathbf{C}$. We assume $\mathfrak{c} \neq \mathfrak{g}$ and take $k$ and $\mu$ as in (i). Given $r \in G_{\mathbf{h}}$, by [S97, Lemma 9.8 (3)] we can put $r \eta_{\mathbf{h}}=\alpha^{-1} t w$ with $\alpha \in G, w \in D\left[\mathfrak{b} \mathfrak{c}, \mathfrak{b}^{-1}\right]$, and $t=\operatorname{diag}\left[q_{1}, \widehat{q}_{1}\right], q_{1} \in G L_{n}(K)_{\mathbf{h}}$. Then, by [S97, Lemma 18.7 (1), (2)] we have

$$
\begin{align*}
& E_{r}(z, s)=\chi_{\mathfrak{c}}\left(\operatorname{det}\left(a_{w}\right)\right)^{-1} j_{\alpha}^{k}(z)^{-1} E_{t}^{*}(\alpha z, s),  \tag{17.26a}\\
& E_{t}^{*}(z, s)=\sum_{h \in S} \operatorname{det}(y)^{-\mu \mathbf{a} / 2} c(h, q, s) \mathbf{e}_{\mathbf{a}}^{n}(h x), \tag{17.26b}
\end{align*}
$$

where $q_{\mathbf{h}}=q_{1}$ and $q_{\mathrm{a}}=y^{1 / 2}$. Now $c(h, q, s)$ is not much different from $c_{h}(y, s)$ of (17.9) in Case SU. Indeed, using the symbols of [S97, Proposition 18.14], put

$$
X=c(S) N(\mathfrak{b c})^{-n^{2}}|\operatorname{det}(y)|^{s \mathbf{a}} \Xi(y, h ;(s+\mu / 2) \mathbf{a},(s-\mu / 2) \mathbf{a})
$$

Then that proposition says that $c(h, q, s) \neq 0$ only if $\left(q^{*} h q\right)_{v} \in(\mathfrak{b c d})_{v}^{-1} \widetilde{S}_{v}$ for every $v \in \mathbf{h}$, in which case

$$
c(h, q, s)=\chi\left(\operatorname{det}\left(q_{1}\right)\right)\left|\operatorname{det}\left(q_{1}\right)\right|_{K}^{n-s} \alpha_{\mathfrak{c}}^{0}\left(\omega q^{*} h q, 2 s, \chi_{1}\right) X
$$

where $\omega$ is an element of $F_{\mathbf{h}}^{\times}$such that $\omega \mathfrak{g}=\mathfrak{b d}$, and $c_{h}(y, s)=\operatorname{det}(y)^{-\mu \mathbf{a} / 2}$ $\cdot c\left(h, y^{1 / 2}, s\right)$, where $c\left(h, y^{1 / 2}, s\right)$ is a special case of $c(h, q, s)$ with $q_{1}=1$. Notice that the quantity $X$ stays the same. Now $\alpha_{\mathrm{c}}^{0}\left(\omega q^{*} h q, 2 s, \chi_{1}\right)$ is given in [S97, Proposition 19.2], and its nature does not depend on $q$. Therefore the analysis of $c(h, q, s)$ at $s=\mu / 2$ is practically the same as that of $c_{h}(y, s)$ in Case SU; the only essential difference is the additional factor $\chi\left(\operatorname{det}\left(q_{1}\right)\right)\left|\operatorname{det}\left(q_{1}\right)\right|_{K}^{n-\mu / 2}$. Combining this with (17.26a, b), we see that

$$
\chi\left(\operatorname{det}\left(q_{1}^{-1}\right) \operatorname{det}\left(a_{w}\right)_{c}\right)\left|\operatorname{det}\left(q_{1}\right)\right|_{K}^{\mu / 2-n} E_{r}(z, \mu / 2)
$$

is of the same nature as $E\left(z, \mu / 2 ; \mu \mathbf{a}, \chi_{1}, \mathfrak{c}\right)$ in Case SU as stated in Theorem 17.7, (iii), (iv), except that we have to consider $\|_{k} \alpha$ with $\alpha$ in $G$, not necessarily in $G_{1}$. Therefore by Theorem 7.11 we obtain (i), (ii), and (iii) when $\mathfrak{c} \neq \mathfrak{g}$. To deal with the case $\mathfrak{c}=\mathfrak{g}$, we use the equality

$$
E_{r}(z, s ; k, \chi, \mathfrak{g})=\sum_{b} E_{r b}(z, s ; k, \chi, \mathfrak{p})
$$

of $\left[\right.$ S97, (19.6.1)], where $\mathfrak{p}$ is a prime ideal and $b$ runs over a finite subset of $G_{\mathbf{h}}$. Thus the results for $\mathfrak{c}=\mathfrak{g}$ follow from those for $\mathfrak{c} \neq \mathfrak{g}$.

Given $k$ and $\mu$ as in (iv), put $p=(m-\mu \mathbf{a}) / 2$ and define $k^{\prime} \in \mathbf{Z}^{\mathbf{b}}$ by $k_{v \rho}^{\prime}=$ $k_{v \rho}-p_{v}$ and $k_{v}^{\prime}=k_{v}-p_{v}$ for $v \in \mathbf{a}$; define also $\Delta_{\mu \mathbf{a}}^{p}$ as in the proof of Theorem 17.9. Then clearly $k_{v \rho}^{\prime}+k_{v}^{\prime}=\mu$ for every $v \in \mathbf{a}$ and moreover

$$
\begin{equation*}
\Delta_{\mu \mathbf{a}}^{p} E_{r}\left(z, \mu / 2 ; k^{\prime}, \chi, \mathfrak{c}\right)=(i / 2)^{n|p|} c_{\mu \mathbf{a}}^{p}(\mu / 2) E_{r}(z, \mu / 2 ; k, \chi, \mathfrak{c}) \tag{17.27}
\end{equation*}
$$

To show this, we first observe that $\Delta_{\mu \mathbf{a}}^{p}\left(f \|_{k^{\prime}} \alpha\right)=\left(\Delta_{\mu \mathbf{a}}^{p} f\right) \|_{k} \alpha$ for every $\alpha \in G$. Indeed, if $\alpha \in\left(G_{1}\right)_{\mathbf{a}}$, this is a special case of (12.21), since $\left\|_{k} \alpha=\right\|_{m} \alpha$ and $\left\|_{k^{\prime}} \alpha=\right\|_{\mu \mathbf{a}} \alpha$. If $\alpha \in G_{\mathbf{a}}$, then $\alpha=c \beta$ with $\beta \in\left(G_{1}\right)_{\mathbf{a}}$ and $c \in \mathbf{C}^{\mathbf{a}}$ such that $\left|c_{v}\right|=1$ for every $v \in \mathbf{a}$, and the equality for $\alpha$ can be reduced to that for $\beta$ (see also §14.4). Now, given $r$, by [S97, Lemma 9.8 (3)] we can put $r=\alpha^{-1} f u$ with $\alpha \in G, f=\operatorname{diag}[\widehat{g}, g], g \in G L_{n}(K)_{\mathbf{h}}$, and $u \in D\left[\mathfrak{b c}, \mathfrak{b}^{-1}\right]$. Then $\alpha_{\mathbf{a}}=u_{\mathbf{a}}$. For $z=x(\mathbf{i})$ with $x \in G_{\mathbf{a}}$ we have

$$
\begin{aligned}
& E_{r}(z, s) j_{x}^{k}(\mathbf{i})^{-1}=E_{\mathbf{A}}(r x, s)=E_{\mathbf{A}}(f u x, s)=\chi_{\mathfrak{c}}\left(\operatorname{det}\left(a_{u}\right)\right)^{-1} E_{\mathbf{A}}\left(f u_{\mathbf{a}} x, s\right) \\
= & \chi_{\mathfrak{c}}\left(\operatorname{det}\left(a_{u}\right)\right)^{-1} E_{f}(\alpha x(\mathbf{i}), s) j_{\alpha x}^{k}(\mathbf{i})^{-1}=\chi_{\mathfrak{c}}\left(\operatorname{det}\left(a_{u}\right)\right)^{-1} E_{f}(\alpha z, s) j_{\alpha}^{k}(z)^{-1} j_{x}^{k}(\mathbf{i})^{-1}
\end{aligned}
$$

Thus we obtain

$$
\begin{equation*}
E_{r}(z, s)=\chi_{\mathfrak{c}}\left(\operatorname{det}\left(a_{u}\right)\right)^{-1} E_{f}(\alpha z, s) j_{\alpha}^{k}(z)^{-1} \tag{17.28a}
\end{equation*}
$$

In [S97, Lemma 18.7 (3)] we showed that

$$
\begin{equation*}
E_{f}(z, s)=\chi\left(\operatorname{det}(g)^{-1}\right)|\operatorname{det}(g)|_{K}^{-s} \sum_{\alpha \in A^{\prime}} N\left(\mathfrak{a}_{f}(\alpha)\right)^{s} \chi[\alpha]_{f} \delta(z)^{s u-m / 2} \|_{k} \alpha \tag{17.28b}
\end{equation*}
$$

with $A^{\prime}=P \backslash\left(G \cap P_{\mathbf{A}} f D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right] f^{-1}\right)$ and $\mathfrak{a}_{f}(\alpha), \chi[\alpha]_{f}$ as described there. These symbols are independent of $k$. Therefore term-wise application of the differential operator shows that (17.27) is true for $E_{f}$ in place of $E_{r}$ for the same reason as in (17.20). Then from (17.28a) we obtain (17.27). Now we have

$$
\begin{equation*}
\Delta_{\mu \mathbf{a}}^{p} \mathcal{N}_{k^{\prime}}^{t}(W) \subset \pi^{n|p|} \mathcal{N}_{k}^{t+n p}(W) \tag{17.29}
\end{equation*}
$$

for a subfield $W$ of $\mathbf{C}$ containing $\Phi$. Therefore we obtain (iv) from (i), (ii), and (17.27), since $c_{\mu \mathbf{a}}^{p}(\mu / 2) \neq 0$ as observed in the proof of Theorem 17.9. Finally, given $\mu$ as in (v), we can repeat the above proof with $D_{r}$ in place of $E_{r}$ if $\mu \geq n$. If $\mu<n$, we put $\nu=2 n-\mu, p=(m-\nu \mathbf{a}) / 2$, and $k^{\prime \prime}=\left(k_{v \rho}-p_{v}, k_{v}-p_{v}\right)_{v \in \mathbf{a}}$. Then $k_{v}^{\prime \prime}+k_{v \rho}^{\prime \prime}=\nu$ for every $v \in \mathbf{a}$, and

$$
\begin{equation*}
\Delta_{\nu \mathbf{a}}^{p} D_{r}\left(z, \mu / 2 ; k^{\prime \prime}, \chi, \mathfrak{c}\right)=(i / 2)^{n|p|} c_{\nu \mathbf{a}}^{p}(\mu / 2) D_{r}(z, \mu / 2 ; k, \chi, \mathfrak{c}) \tag{17.30}
\end{equation*}
$$

Thus we obtain (v) for the same reason as above.
17.13. Lemma. Define $E_{r}(z, s)$ by (17.23a) for $r \in G_{\mathbf{h}}$ in Cases $S P$ and UT with $k \in \mathbf{Z}^{\mathbf{b}}$ and $\kappa=0$. Then there is a finite sum expression

$$
\begin{equation*}
E_{r}(z, s)=\sum_{i=1}^{m} b_{i} c_{i}^{s} E\left(z, s ; k, \Gamma_{i}\right) \|_{k} \alpha_{i} \tag{17.31}
\end{equation*}
$$

with $b_{i} \in \overline{\mathbf{Q}}, 0<c_{i} \in \mathbf{R} \cap \overline{\mathbf{Q}}$, congruence subgroups $\Gamma_{i}$ of $G$, and $\alpha_{i} \in G$. Moreover, in Case $U T$ we can take $\Gamma_{i}$ in $S U\left(\eta_{n}\right)$.

Proof. In Case UT we have (17.28a, b) with $f=\operatorname{diag}[\hat{g}, g]$. By [S97, Lemma 18.7 (4)], $A^{\prime}$ of (17.28b) can be given by $\bigsqcup_{\beta \in B} S_{\beta} \beta$ with a finite subset $B$ of $G$ and $S_{\beta}=\left(P \cap \beta \Gamma \beta^{-1}\right) \backslash \beta \Gamma \beta^{-1}$, where $\Gamma=G \cap f D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right] f^{-1}$. For $\gamma \in S_{\beta}$ we have $\mathfrak{a}_{f}(\gamma \beta)=\mathfrak{a}_{f}(\beta)$ and $\chi[\gamma \beta]_{f}=c_{\beta} \chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{\gamma}\right)\right)$ with a constant $c_{\beta}$ independent of $\gamma$. Therefore, taking a suitable congruence subgroup of $\beta \Gamma \beta^{-1}$, we obtain expression (17.31) for $E_{f}$, which combined with (17.28a) proves (17.31) in the general case. Now (17.5) is valid in Case UT, and in that formula we can take $\Gamma_{1}$ in $S U\left(\eta_{n}\right)$. This proves our lemma in Case UT. Case SP is simpler, since we have $G_{\mathbf{A}}=G D\left[\mathfrak{b}^{-1}, \mathfrak{b}\right]$, so that (17.28a) holds with $f=1$, and $E_{1}(z, s)$ is given by (16.39).

## 18. Eisenstein series in the Hilbert modular case

18.1. In this section we consider Eisenstein series when the group $G$ is $S L_{2}(F)$, and as applications we treat critical values of certain $L$-functions of a totally imaginary quadratic extension of $F$. Though our results on Eisenstein series are essentially special cases of what was obtained in Sections 16 and 17, we present them in a somewhat different fashion. We start with some elementary facts on $L$-functions of $F$ stated in the form suitable for our later applications. Throughout this section we denote by $\mathfrak{d}, D_{F}$, and $R_{F}$ the different of $F$ relative to $\mathbf{Q}$, the discriminant of $F$, and the regulator of $F$. We put $[F: \mathbf{Q}]=e$ and $N(x)=N_{F / \mathbf{Q}}(x)$ for $x \in F$.

Then $N(x \mathfrak{g})=|N(x)|=|x|^{\mathbf{a}}$ for $x \in F^{\times}$. For $a_{0} \in F$ and a $\mathfrak{g}$-ideal $\mathfrak{a}$ we put $a_{0}+\mathfrak{a}=\left\{a_{0}+a \mid a \in \mathfrak{a}\right\}$.

Given $\kappa \in \mathcal{S}\left(F_{\mathbf{h}}\right)$ and $k \in \mathbf{Z}^{\mathbf{a}}$, we can find a subgroup $U$ of $\mathfrak{g}^{\times}$of finite index such that $\kappa(u d)=\kappa(d)$ and $u^{-k}|u|^{k}=1$ for every $u \in U$ and every $d \in F$. As explained in $\S 1.6$, we view $\kappa$ as a function on $F$. With such a $U$ we put

$$
\begin{equation*}
D_{k}(s, \kappa)=\left[\mathfrak{g}^{\times}: U\right]^{-1} \sum_{d \in F^{\times} / U} \kappa(d) d^{-k}|d|^{k-s \mathbf{a}} \quad\left(s \in \mathbf{C}^{\times}\right) . \tag{18.1}
\end{equation*}
$$

Clearly this does not depend on the choice of $U$, and is convergent for $\operatorname{Re}(s)>1$.
18.2. Lemma. Let $t_{v}=0$ or 1 according as $k_{v}$ is even or odd. Put $R_{t}(s, \kappa)$ $=\Gamma_{t}(s) D_{k}(s, \kappa)$ with $\Gamma_{t}(s)=\prod_{v \in \mathbf{a}} \pi^{-\left(s+t_{v}\right) / 2} \Gamma\left(\left(s+t_{v}\right) / 2\right)$. Then $R_{t}(s, \kappa)$ can be continued as a meromorphic function of $s$ to the whole $\mathbf{C}$, and satisfies $R_{t}(1-s, \kappa)=R_{t}\left(s, \kappa_{*}\right)$ with the element $\kappa_{*}$ of $\mathcal{S}\left(F_{\mathbf{h}}\right)$ given by

$$
\kappa_{*}(x)=i^{-|t|} \int_{F_{\mathbf{h}}} \kappa(y) \mathbf{e}_{\mathbf{h}}(-x y) d y
$$

where $|t|=\sum_{v \in \mathbf{a}} t_{v}$ and dy is the Haar measure of $F_{\mathbf{h}}$ such that $\prod_{v \in \mathbf{h}} \mathfrak{g}_{v}$ has measure $D_{F}^{-1 / 2}$. Moreover, $R_{t}(s, \kappa)$ is entire except when $t=0$, in which case $R_{t}(s, \kappa)$ is holomorphic on $\mathbf{C}$ except for possible simple poles at $s=0$ and $s=1$ with residues $-2^{e-1} \kappa(0) R_{F}$ and $2^{e-1} \kappa_{*}(0) R_{F}$, respectively.

Though this can be proved by the standard method by taking a suitable zeta integral on $F_{\mathbf{A}}^{\times}$, we shall derive it from a more general result concerning the Mellin transform of a Hilbert modular form in §A7.3.

Take $\kappa$ to be the characteristic function of the set $a_{0}+\mathfrak{a}$ as above, for example. Then $\kappa_{*}(0)=i^{-|t|} D_{F}^{-1 / 2} N(\mathfrak{a})^{-1}$, and so if $t=0$, (that is, if $k \in 2 \mathbf{Z}^{\mathbf{a}}$,) then $D_{k}(s, \kappa)$ has a simple pole at $s=1$ with residue $2^{e-1} D_{F}^{-1 / 2} N(\mathfrak{a})^{-1} R_{F}$.
18.3. We put $P=\left\{\alpha \in G \mid c_{\alpha}=0\right\}$ and $\mathcal{H}=\mathfrak{H}_{1}^{\text {a }}$ as before. In addition, we put $H=F_{2}^{1}-\{0\}$ and

$$
\begin{equation*}
j_{h}(z)=\left(c_{v} z_{v}+d_{v}\right)_{v \in \mathbf{a}} \text { for } h=(c, d) \in H \text { and } z \in \mathcal{H} \tag{18.2}
\end{equation*}
$$

We easily see that $j_{h \gamma}(z)=j_{h}(\gamma z) j_{\gamma}(z)$ for $h \in H$ and $\gamma \in G$. Given $k \in \mathbf{Z}^{\mathbf{a}}$ and $\lambda \in \mathcal{S}\left(\left(F_{2}^{1}\right)_{\mathbf{h}}\right)$, we can find a subgroup $U$ of $\mathfrak{g}^{\times}$of finite index such that $u^{k}|u|^{-k}=1$ and $\lambda(u h)=\lambda(h)$ for every $u \in U$ and every $h \in F_{2}^{1}$. Then we put, for $(z, s) \in \mathcal{H} \times \mathbf{C}$,

$$
\begin{align*}
E_{k}(z, s ; \lambda) & =\left[\mathfrak{g}^{\times}: U\right]^{-1} \sum_{h \in H / U} \lambda(h) j_{h}(z)^{-k} y^{s \mathbf{a}-k / 2}\left|j_{h}(z)\right|^{k-2 s \mathbf{a}}  \tag{18.3}\\
& =\left[\mathfrak{g}^{\times}: U\right]^{-1} \sum_{(c, d) \in H / U} \lambda(c, d)(c z+d)^{-k} y^{s \mathbf{a}-k / 2}|c z+d|^{k-2 s \mathbf{a}},
\end{align*}
$$

where of course $y=\operatorname{Im}(z)$. The sum is formally well-defined, independently of the choice of $U$. Defining $\lambda^{\gamma}$ for $\gamma \in G$ by $\lambda^{\gamma}(h)=\lambda\left(h \gamma^{-1}\right)$, we can easily verify that .

$$
\begin{equation*}
E_{k}\left(z, s ; \lambda^{\gamma}\right)=E_{k}(\gamma z, s ; \lambda) j_{\gamma}(z)^{-k} \text { for every } \gamma \in G \tag{18.4}
\end{equation*}
$$

If we put $\Gamma=\left\{\gamma \in G \mid \lambda^{\gamma}=\lambda\right\}$, then clearly $\Gamma$ is a congruence subgroup of $G$, and $E_{k}(z, s ; \lambda) \|_{k} \gamma=E_{k}(z, s ; \lambda)$ for every $\gamma \in \Gamma$.

Let us now show that $E_{k}(z, s ; \lambda)$ is a finite "linear combination" of transforms of functions of type (17.3). We first put $x_{0}=\left[\begin{array}{ll}0 & 1\end{array}\right]$ and observe that $x_{0} \alpha=\left[\begin{array}{ll}c_{\alpha} & d_{\alpha}\end{array}\right]$
for $\alpha \in G, P=\left\{\alpha \in G \mid x_{0} \alpha \in F x_{0}\right\}$, and the map $\alpha \mapsto x_{0} \alpha$ gives a bijection of $P \backslash G$ onto $F^{\times} \backslash H$. Then we easily see that the map $(d, \alpha) \mapsto d x_{0} \alpha$ for $(d, \alpha) \in$ $F^{\times} \times G$ gives a bijection of $\left(F^{\times} / U\right) \times(P \backslash G)$ onto $H / U$.

Given $\lambda \in \mathcal{S}\left(\left(F_{2}^{1}\right)_{\mathbf{h}}\right)$, take $U$ and $\Gamma$ as above, and take also a complete set of representatives $B$ for $P \backslash G / \Gamma$, which is finite by [S97, Lemma 9.8 (3)]. Then $P \backslash G$ can be given by $\bigsqcup_{\beta \in B} \mathcal{R}_{\beta} \beta$ with $\mathcal{R}_{\beta}=\left(P \cap \beta \Gamma \beta^{-1}\right) \backslash \beta \Gamma \beta^{-1}$. Write an element $h$ of $H / U$ in the form $h=d x_{0} \gamma \beta$ with $d \in F^{\times} / U, \beta \in B$, and $\gamma \in \mathcal{R}_{\beta}$. Since $\lambda\left(d x_{0} \gamma \beta\right)=\lambda\left(d x_{0} \beta \beta^{-1} \gamma \beta\right)=\lambda\left(d x_{0} \beta\right)$, we have

$$
\begin{align*}
E_{k}(z, s ; \lambda) & =\left[\mathfrak{g}^{\times}: U\right]^{-1} \sum_{\beta \in B} \sum_{d \in F^{\times} / U} \lambda\left(d x_{0} \beta\right) d^{-k}|d|^{k-2 s \mathbf{a}} \sum_{\gamma \in \mathcal{R}_{\beta}} y^{s \mathbf{a}-k / 2} \|_{k} \gamma \beta  \tag{18.5}\\
& =\sum_{\beta \in B} D_{k}\left(2 s, \kappa_{\beta}\right) E\left(z, s ; k, \beta \Gamma \beta^{-1}\right) \|_{k} \beta
\end{align*}
$$

where $\kappa_{\beta}$ is defined by $\kappa_{\beta}(d)=\lambda\left(d x_{0} \beta\right)$ for $d \in F$ and $E\left(z, s ; k, \beta \Gamma \beta^{-1}\right)$ by (17.3a). Since (17.3) is convergent for $\operatorname{Re}(s)>1$ and $D_{k}\left(2 s, \kappa_{\beta}\right)$ is holomorphic for $\operatorname{Re}(s)>1 / 2$, we see that $E_{k}(z, s ; \lambda)$ is meaningful as a holomorphic function of $s$ at least for $\operatorname{Re}(s)>1$. Also, from Lemmas 17.2 and 18.2 we see that $E_{k}(z, s ; \lambda)$ can be continued as a meromorphic function of $s$ to the whole $\mathbf{C}$. We can actually derive a stronger result from the explicit Fourier expansion of $E_{k}(z, s ; \lambda)$, which is our next problem. For that purpose we need
18.4. Lemma. For a $\mathfrak{g}$-ideal $\mathfrak{m}$ we have

$$
D_{F}^{1 / 2} N(\mathfrak{m}) \sum_{a \in \mathfrak{m}}(z+a)^{-\alpha}(\bar{z}+a)^{-\beta}=\sum_{b \in(\mathfrak{o} \mathfrak{m})^{-1}} \mathbf{e}_{\mathbf{a}}(b x) \Xi(y, b ; \alpha, \beta) .
$$

Here $z=x+i y \in \mathcal{H}, \alpha \in \mathbf{C}^{\mathbf{a}}, \beta \in \mathbf{C}^{\mathbf{a}}, \operatorname{Re}\left(\alpha_{v}+\beta_{v}\right)>1$ for every $v \in \mathbf{a}$, and $\Xi(y, b ; \alpha, \beta)=\prod_{v \in \mathbf{a}} \xi\left(y_{v}, b_{v} ; \alpha_{v}, \beta_{v}\right)$ with

$$
\xi\left(g, h ; s, s^{\prime}\right)=\int_{\mathbf{R}} \mathbf{e}(-h t)(t+i g)^{-s}(t-i g)^{-s^{\prime}} d t \quad\left(s, s^{\prime} \in \mathbf{C} ; 0<g \in \mathbf{R}, h \in \mathbf{R}\right)
$$

Proof. Put $f(x)=(x+i y)^{-\alpha}(x-i y)^{-\beta}$ for $x \in \mathbf{R}^{\mathbf{a}}$ with a fixed $y \in \mathbf{R}^{\mathbf{a}}, \gg 0$, and let $\hat{f}(x)$ be its Fourier transform. (That $f$ is $L^{1}$ will be shown later.) Then $\hat{f}(x)=\Xi(y, x ; \alpha, \beta)$. By the Poisson summation formula we obtain $\operatorname{vol}\left(F_{\mathbf{a}} / \mathfrak{m}\right)$ $\cdot \sum_{a \in \mathfrak{m}} f(x+a)=\sum_{b \in(\mathfrak{o m})^{-1}} \mathbf{e}_{\mathbf{a}}(b x) \hat{f}(b)$, which gives the desired equality, and which holds if the left-hand side is convergent and defines a $C^{\infty}$ function of $x$. To see the last point, we take an easy equality

$$
|x+i y|^{-2 r}=\prod_{v \in \mathbf{a}} \Gamma\left(r_{v}\right)^{-1} \int_{0}^{\infty} e^{-t\left|z_{v}\right|^{2}} t^{r_{v}-1} d t
$$

valid for $z=x+i y \in \mathbf{C}^{\mathbf{a}}$ and $r \in \mathbf{R}^{\mathbf{a}}, \gg 0$. Then we can easily verify that

$$
\int_{\mathbf{R}^{\mathbf{a}}}|x+i y|^{-2 r} d x=\pi^{e / 2} y^{\mathbf{a}-2 r} \prod_{v \in \mathbf{a}} \Gamma\left(r_{v}-1 / 2\right) \Gamma\left(r_{v}\right)^{-1}
$$

if $2 r-\mathbf{a} \gg 0$. Let $T$ be a compact subset of $\mathcal{H}$. By a well-known principle, there exists an open subset $U$ of $\mathcal{H}$ containing $T$, whose closure is compact and contained in $\mathcal{H}$, and a constant $C_{1}$ such that

$$
\begin{equation*}
|h(w)| \leq C_{1} \int_{U}|h(z)| d v(z) \tag{*}
\end{equation*}
$$

for every $w \in T$ and every holomorphic function $h$ on $\mathcal{H}$, where $d v(z)$ is the Lebesgue measure on $\mathbf{C}^{\mathbf{a}}$. Take a compact subset $X$ of $R^{\mathbf{a}}$ and a compact subset $Y$
of $\left\{y \in \mathbf{R}^{\mathbf{a}} \mid y \gg 0\right\}$ so that $U \subset X+i Y$. Let $M=\#\{a \in \mathfrak{m} \mid(a+X) \cap X \neq \varnothing\}$. Then, for any finite subset $A$ of $\mathfrak{m}$, we have
$\int_{Y} \int_{X_{a \in A}} \sum_{a}|x+i y+a|^{-2 r} d x d y \leq M \int_{Y} \int_{\mathbf{R}^{\mathbf{a}}}|x+i y|^{-2 r} d x d y \leq M C_{2} \int_{Y} y^{\mathbf{a}-2 r} d y \leq M C_{3}$ if $r$ belongs to a compact subset $Z$ of $\left\{r \in \mathbf{R}^{\mathbf{a}} \mid 2 r-\mathbf{a} \gg 0\right\}$, where $C_{2}$ and $C_{3}$ are constants depending only on $Y$ and $Z$. By ( ${ }^{*}$ ) we have $\sum_{a \in A}|z+a|^{-2 r} \leq C_{4}$ for every $z \in T$ and every $r \in Z$ with a constant $C_{4}$ depending only on $X, Y$, and $Z$. Now $\left|w^{-\alpha} \bar{w}^{-\beta}\right| \leq C_{5}|w|^{-\operatorname{Re}(\alpha+\beta)}$ for every $w \in \mathcal{H}$ if $(\alpha, \beta)$ belongs to a compact subset $S$ of $\mathbf{C}^{\mathbf{a}} \times \mathbf{C}^{\mathbf{a}}$, where $C_{5}$ is a constant depending only on $S$. Thus $f$ belongs to $L^{1}\left(\mathbf{R}^{\mathbf{a}}\right)$ and $\sum_{a \in \mathrm{~m}} f(x+a)$ is locally uniformly convergent for ( $x, a, \beta$ ) if $\operatorname{Re}(\alpha+\beta) \gg \mathbf{a}$. Since $\left(\partial / \partial x_{v}\right)\left(z^{-\alpha} \bar{z}^{-\beta}\right)=-\alpha_{v} z^{-\alpha-v} \bar{z}^{-\beta}-\beta_{v} z^{-\alpha} \bar{z}^{-\beta-v}$, we see that $\sum_{a \in \mathfrak{m}}\left(\partial / \partial x_{v}\right) f(x+a)$ is locally uniformly convergent if $\operatorname{Re}(\alpha+\beta) \gg \mathbf{a}$. The same is true for derivatives of any order. This proves that $\sum_{a \in \mathfrak{m}} f(x+a)$ converges to a $C^{\infty}$ function if $\operatorname{Re}(\alpha+\beta) \gg \mathbf{a}$, which completes the proof.

Notice that the above $\xi$ is exactly the function of (16.44) with $n=1$. We insert here an easy fact:
18.5. Lemma. Let $f(z)=\sum_{h \in F} b(h) \mathbf{e}_{\mathbf{a}}(h z)$ and $g(z)=\sum_{h \in F} c(h) \mathbf{e}_{\mathbf{a}}(h z)$ be elements of $\mathcal{M}_{\mu \mathbf{a}}$ with $0<\mu \in 2^{-1} \mathbf{Z}$, and let $\sigma \in \operatorname{Aut}(\mathbf{C})$. If $b(h)^{\sigma}=c(h)$ for every $h \in F, \neq 0$, then $b(0)^{\sigma}=c(0)$. Consequently $b(0)$ is contained in the field generated over $\mathbf{Q}$ by the $b(h)$ for all $h \neq 0$.

Proof. By Theorem 9.9 (4) or Theorem 10.7 (5), $f^{\sigma} \in \mathcal{M}_{\mu \mathbf{a}}$, and hence $b(0)^{\sigma}-$ $c(0)=f^{\sigma}-g \in \mathcal{M}_{\mu \mathbf{a}}$. Since $\mu>0$, we have $b(0)^{\sigma}-c(0)=0$, which proves our lemma.
18.6. We now take $\lambda$ in (18.3) to be the characteristic function of $\left(a_{0}+\mathfrak{a}\right) \times$ $\left(b_{0}+\mathfrak{b}\right)$ with $a_{0}, b_{0} \in F$ and $\mathfrak{g}$-ideals $\mathfrak{a}$ and $\mathfrak{b}$. Then we can take

$$
U=\left\{u \in \mathfrak{g}^{\times} \mid u \gg 0,(u-1) a_{0} \in \mathfrak{a},(u-1) b_{0} \in \mathfrak{b}\right\} .
$$

Notice that $\mathcal{S}\left(\left(F_{2}^{1}\right)_{\mathbf{h}}\right)$ is spanned by such $\lambda$ 's. Let $A=\left[\left(a_{0}+\mathfrak{a}\right) \cap F^{\times}\right] / U$ and $B=\left[\left(b_{0}+\mathfrak{b}\right) \cap F^{\times}\right] / U$; also let $\kappa_{A}$ resp. $\kappa_{B}$ be the characteristic function of $a_{0}+\mathfrak{a}$ resp. $b_{0}+\mathfrak{b}$. Then

$$
\begin{aligned}
{\left[\mathfrak{g}^{\times}: U\right] y^{k / 2-s \mathbf{a}} E_{k}(z, s ; \lambda)=} & \varepsilon\left(a_{0}, \mathfrak{a}\right) \sum_{d \in B} d^{-k}|d|^{k-2 s \mathbf{s}} \\
& +\sum_{c \in A} \sum_{d \in b_{0}+\mathfrak{b}}(c z+d)^{-k}|c z+d|^{k-2 s \mathbf{a}}
\end{aligned}
$$

where $\varepsilon\left(a_{0}, \mathfrak{a}\right)$ is 1 or 0 according as $a_{0} \in \mathfrak{a}$ or $a_{0} \notin \mathfrak{a}$. For a fixed $c \in A$ we have

$$
\sum_{d \in b_{0}+\mathfrak{b}}(c z+d)^{-k}|c z+d|^{k-2 s \mathbf{a}}=c^{-k}|c|^{k-2 s \mathbf{a}} \sum_{a \in c^{-1} \mathfrak{b}}\left(z+c^{-1} b_{0}+a\right)^{-k}\left|z+c^{-1} b_{0}+a\right|^{k-2 s \mathbf{a}} .
$$

By Lemma 18.4 this equals

$$
c^{-k}|c|^{k-2 s \mathbf{a}} D_{F}^{-1 / 2} N\left(c \mathfrak{b}^{-1}\right) \sum_{b \in c \mathfrak{b}^{-1} \mathfrak{d}^{-1}} \mathbf{e}_{\mathbf{a}}\left(b x+b c^{-1} b_{0}\right) \Xi(y, b ; s \mathbf{a}+k / 2, s \mathbf{a}-k / 2)
$$

Thus the Fourier expansion of $E_{k}$ can be given by

$$
\begin{align*}
& \text { 6) } \quad E_{k}(z, s ; \lambda)=\varepsilon\left(a_{0}, \mathfrak{a}\right) y^{\mathbf{s a}-k / 2} D_{k}\left(2 s, \kappa_{B}\right)  \tag{18.6}\\
& +D_{F}^{-1 / 2} N(\mathfrak{b})^{-1} y^{\mathbf{s a}-k / 2} \Xi(y, 0 ; s \mathbf{a}+k / 2, s \mathbf{a}-k / 2) D_{k}\left(2 s-1, \kappa_{A}\right)
\end{align*}
$$

$$
\begin{gathered}
+\left[\mathfrak{g}^{\times}: U\right]^{-1} D_{F}^{-1 / 2} N(\mathfrak{b})^{-1} y^{s \mathbf{a}-k / 2} \sum_{b \in F^{\times}} \mathbf{e}_{\mathbf{a}}(b x) \Xi(y, b ; s \mathbf{a}+k / 2, s \mathbf{a}-k / 2) \varphi(b, s) \\
\text { with } \varphi(b, s)=\sum_{a, c} \mathbf{e}_{\mathbf{a}}\left(a b_{0}\right) c^{-k}|c|^{k+(1-2 s) \mathbf{a}}
\end{gathered}
$$

where the last sum is taken over all $(a, c) \in \mathfrak{b}^{-1} \mathfrak{d}^{-1} \times A$ such that $a c=b$. Notice that it is a finite sum.
18.7. Proposition. The product $E_{k}(z, s ; \lambda) \prod_{v \in \mathbf{a}} \Gamma\left(s+\left|k_{v}\right| / 2\right)$ can be continued as a meromorphic function of $s$ to the whole $\mathbf{C}$, which is holomorphic except for possible simple poles at $s=0$ and $s=1$. The pole at $s=1$ occurs if and only if $k=0$, and the residue is $2^{e-2} \pi^{e} N(\mathfrak{a b d})^{-1} R_{F}$; the pole at $s=0$ occurs if and only if $k=0, a_{0} \in \mathfrak{a}$, and $b_{0} \in \mathfrak{b}$, and the residue is $-2^{e-2} R_{F}$.

Proof. Put $\Delta(s)=\prod_{v \in \mathbf{a}} \Gamma\left(s+\left|k_{v}\right| / 2\right)$ and $E_{k}(z, s ; \lambda)=\sum_{b \in F} c_{b}(y, s) \mathbf{e}_{\mathbf{a}}(b x)$. We already know meromorphic continuation of $E_{k}(z, s ; \lambda)$ to the whole $\mathbf{C}$ (which can actually be derived from the above Fourier expansion). Thus our task is to study $\Delta(s) c_{b}(y, s)$ for each $b$. From (17.11) and the above formula for $\varphi(b, s)$, we see that $\Delta(s) c_{b}(y, s)$ is holomorphic everywhere if $b \neq 0$. Now, from (17.13) we obtain

$$
\begin{equation*}
\xi\left(g, 0 ; s+k_{v} / 2 ; s-k_{v} / 2\right)=i^{-k_{v}} 2^{2-2 s} \pi g^{1-2 s} \frac{\Gamma(2 s-1)}{\Gamma\left(s+k_{v} / 2\right) \Gamma\left(s-k_{v} / 2\right)} \tag{18.7}
\end{equation*}
$$

For $b=0, \Delta(s) c_{0}(y, s)$ consists of two terms: one involving $\Delta(s) D_{k}\left(2 s, \kappa_{B}\right)$ and the other involving $\Delta(s) D_{k}\left(2 s-1, \kappa_{A}\right)$. From Lemma 18.2 we easily see that the former has possible simple poles at $s=0$ and $s=1 / 2$; similarly, by (18.7), the latter has possible simple poles at $s=1 / 2$ and $s=1$. Therefore $\Delta(s) E_{k}(z, s ; \lambda)$ may have poles only at these points, and the residue at $s=0,1 / 2$, or 1 is a constant times $y^{-k / 2}, y^{(\mathbf{a}-k) / 2}$, or $y^{-k / 2}$, respectively. However, by (18.4), these must be invariant under the action of $\|_{k} \gamma$ for every $\gamma \in \Gamma$, which is possible only when $s \neq 1 / 2$ and $k=0$, and the residue is a constant. Therefore $\Delta(s) E_{k}(z, s ; \lambda)$ is entire if $k \neq 0$. If $k=0$, it has possible simple poles at $s=0$ and $s=1$. The residue at $s=0$ is that of $\varepsilon\left(a_{0}, \mathfrak{a}\right) R_{0}\left(2 s, \kappa_{B}\right)$ at $s=0$. By Lemma 18.2 this is nonzero if and only if $a_{0} \in \mathfrak{a}$ and $b_{0} \in \mathfrak{b}$. Similarly, by (18.7) and Lemma 18.2, the residue at $s=1$ is a nonzero constant times the residue of $R_{0}\left(2 s-1, \kappa_{A}\right)$ at $s=1$, which is nonzero. A simple calculation gives each residue as stated in our proposition.
18.8. We now assume that $k=\mu \mathbf{a}$ with $0<\mu \in \mathbf{Z}$, and take the value of $E_{k}$ at $s=\mu / 2$. From (17.12) and (18.7) we obtain

$$
\begin{aligned}
& \xi(g, h ; \mu, 0)= \begin{cases}(-2 \pi i)^{\mu} \Gamma(\mu)^{-1} h^{\mu-1} \mathbf{e}(i g h) & \text { if } h>0 \\
0 & \text { if } h<0\end{cases} \\
& \lim _{\sigma \rightarrow 0} \xi(g, 0 ; \sigma+\mu, \sigma)= \begin{cases}-\pi i & \text { if } \mu=1, \\
0 & \text { if } \mu>1,\end{cases} \\
& \lim _{\sigma \rightarrow 0} \xi(g, 0 ; \sigma+2, \sigma) / \sigma=-\pi g^{-1}
\end{aligned}
$$

Putting $k=\mu \mathbf{a}$ and $s=\mu / 2$ in (18.6), we have

$$
\begin{align*}
& D_{F}^{1 / 2} N(\mathfrak{b})(-2 \pi i)^{-\mu e} \Gamma(\mu)^{e} E_{\mu \mathbf{a}}(z, \mu / 2 ; \lambda)  \tag{18.8}\\
& \quad=\left[\mathfrak{g}^{\times}: U\right]^{-1} \sum_{0 \ll b \in F} \varphi(b) \mathbf{e}_{\mathbf{a}}(b z)+\mathcal{D}_{1}(\mu)+\mathcal{D}_{2}(0)+C \\
& \text { with } \varphi(b)=\sum_{0 \neq a \in \mathfrak{b}^{-1} \mathfrak{d}^{-1}, a^{-1} b \in A} N(a)^{\mu}|N(a)|^{-1} \mathbf{e}_{\mathbf{a}}\left(a b_{0}\right)
\end{align*}
$$

$$
\begin{aligned}
& \mathcal{D}_{1}(s)=\varepsilon\left(a_{0}, \mathfrak{a}\right) D_{F}^{1 / 2} N(\mathfrak{b})(-2 \pi i)^{-\mu e} \Gamma(\mu)^{e} D_{\mu \mathbf{a}}\left(s, \kappa_{B}\right), \\
& \mathcal{D}_{2}(s)= \begin{cases}2^{-e} D_{\mathbf{a}}\left(s, \kappa_{A}\right) & \text { if } \mu=1, \\
0 & \text { if } \mu>1,\end{cases} \\
& C= \begin{cases}(8 \pi y)^{-1} N(\mathfrak{a})^{-1} & \text { if } \quad F=\mathbf{Q} \text { and } \mu=2, \\
0 & \text { otherwise } .\end{cases}
\end{aligned}
$$

Notice that $\Xi(y, 0 ;(s+2) \mathbf{a}, s \mathbf{a})$ has a zero of order $e$ at $s=0$, which is why $C=0$ if $F \neq \mathbf{Q}$.

Next let us consider $E_{\mu \mathbf{a}}(z, \nu / 2 ; \lambda)$ for $\nu=2-\mu$. We note that

$$
\xi(g, h ; 1,1-\mu)= \begin{cases}(-2 i)^{\mu} \pi g^{\mu-1} \mathbf{e}(i g h) & \text { if } h>0 \\ (-2 i)^{\mu} 2^{-1} \pi g^{\mu-1} & \text { if } h=0 \\ 0 & \text { if } h<0\end{cases}
$$

These can be obtained from (17.11) and the formula $\omega(g, h ; 1, \beta)=2^{-1} \mathbf{e}(i g h)$ given in [S82, (4.35K)]. If $\mu=1$, then $\nu=1$ and $E_{\mu \mathbf{a}}(z, \nu / 2 ; \lambda)=E_{\mu \mathbf{a}}(z, \mu / 2 ; \lambda)$, which is already given. Therefore we assume $\mu>1$, so that $\nu \leq 0$. Then

$$
\begin{gather*}
\quad D_{F}^{1 / 2} N(\mathfrak{b})(-2 i)^{-\mu e} \pi^{-e} E_{\mu \mathbf{a}}(z, \nu / 2 ; \lambda)  \tag{18.9}\\
=\left[\mathfrak{g}^{\times}: U\right]^{-1} \sum_{0 \ll b \in F} \varphi^{\prime}(b) \mathbf{e}_{\mathbf{a}}(b z)+2^{-e} D_{\mu \mathbf{a}}\left(\nu-1, \kappa_{A}\right)+C^{\prime} \\
\text { with } \quad \varphi^{\prime}(b)=\sum_{c \in A, c^{-1} b \in \mathfrak{b}^{-1} \mathfrak{d}^{-1}}|N(c)|^{-1} N(c)^{\mu} \mathbf{e}_{\mathbf{a}}\left(c^{-1} b b_{0}\right), \\
C^{\prime}= \begin{cases}\varepsilon\left(a_{0}, \mathfrak{a}\right) \varepsilon\left(b_{0}, \mathfrak{b}\right) N(\mathfrak{b})(4 \pi y)^{-1} & \text { if } F=\mathbf{Q} \text { and } \mu=2, \\
0 & \text { otherwise. }\end{cases}
\end{gather*}
$$

18.9. Theorem. Let $\lambda$ be a $\mathbf{Q}_{\mathrm{ab}}$-valued element of $\mathcal{S}\left(\left(F_{2}^{1}\right)_{\mathbf{h}}\right)$ and let $0<\mu \in \mathbf{Z}$. Then the following assertions hold:
(1) $E_{\mu \mathbf{a}}(z, \mu / 2 ; \lambda)$ belongs to $\pi^{\mu e} \mathcal{M}_{\mu \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ except when $F=\mathbf{Q}$ and $\mu=2$, in which case if belongs to $\pi^{2} \mathcal{N}_{2}^{1}\left(\mathbf{Q}_{\mathrm{ab}}\right)$.
(2) $E_{\mu \mathbf{a}}(z, 1-\mu / 2 ; \lambda)$ belongs to $\pi^{e} \mathcal{M}_{\mu \mathbf{a}}\left(\mathbf{Q}_{\mathbf{a b}}\right)$ except when $F=\mathbf{Q}$ and $\mu=2$, in which case it belongs to $\pi \mathcal{N}_{2}^{1}\left(\mathbf{Q}_{\mathrm{ab}}\right)$.

Proof. We may assume that $\lambda$ is the characteristic function of $\left(a_{0}+\mathfrak{a}\right) \times$ $\left(b_{0}+\mathfrak{b}\right)$. From (18.4) and (18.8) we see that $E_{\mu \mathbf{a}}(z, \mu / 2 ; \lambda)$ belongs to $\mathcal{M}_{\mu \mathbf{a}}$ except when $F=\mathbf{Q}$ and $\mu=2$. Since $\varphi(b) \in \mathbf{Q}_{\mathrm{ab}}$, Lemma 18.5 shows that $\pi^{-\mu e} E_{\mu \mathbf{a}}(z, \mu / 2 ; \lambda) \in \mathcal{M}_{\mu \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$. Suppose $F=\mathbf{Q}$ and $\mu=2$. Then (18.4) and (18.8) show that $E_{2 \mathbf{a}}(z, 1 ; \lambda)$ satisfies condition (13.18a) with $p=1$. Thus $E_{2 \mathbf{a}}(z, 1 ; \lambda) \in$ $\mathcal{N}_{2}^{1}$. Now $\varphi(b) \in \mathbf{Q}_{\mathrm{ab}}$, and from that fact we can conclude that $E_{2 \mathbf{a}}(z, 1 ; \lambda) \in$ $\pi^{2} \mathcal{N}_{2}^{1}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ for the reason which will be explained in the proof of the following proposition. This proves (1). Assertion (2) can be proved in the same manner.

Remark. In (18.3) take $F=\mathbf{Q}, U=\{ \pm 1\}, k=2$, and take $\lambda$ to be the characteristic function of $\prod_{p} \mathbf{Z}_{p}$, which means that $\mathfrak{a}=\mathfrak{b}=\mathbf{Z}$ and $a_{0}=b_{0}=0$. Then $D_{2 \mathbf{a}}\left(s, \kappa_{A}\right)=\zeta(s)$, and (18.8) with $\mu=2$ shows that

$$
(2 \pi i)^{-2} E_{2}(z, 1 / 2 ; \lambda)=(8 \pi y)^{-1}-(1 / 24)+\sum_{b=1}^{\infty} \mathbf{e}_{\mathbf{a}}(b z) \sum_{0<a \mid b} a
$$

which is exactly twice formula (0.7) of the introduction. This type of nearly holomorphic Eisenstein series occurs in Cases SP and UT if $F=\mathbf{Q}$, as shown in Theorem
17.7 (iv). For further investigations of the series of this nature, see $[\mathrm{S} 83, \S 9]$ and [S85a, pp.290-291].
18.10. Proposition. For $0<\mu \in \mathbf{Z}, b_{0} \in F$, and $a \mathfrak{g}$-ideal $\mathfrak{b}$ put $D_{\mu}\left(s ; b_{0}, \mathfrak{b}\right)=$ $D_{\mu \mathbf{a}}(s, \kappa)$, where $\kappa$ is the characteristic function of $b_{0}+\mathfrak{b}$; put also $Q\left(\mu ; b_{0}, \mathfrak{b}\right)=$ $(2 \pi i)^{-\mu e} D_{F}^{1 / 2} D_{\mu}\left(\mu ; b_{0}, \mathfrak{b}\right)$. Then the following assertions hold:
(1) $Q\left(\mu ; b_{0}, \mathfrak{b}\right) \in \mathbf{Q}_{\mathrm{ab}}$. Moreover, let $\sigma=[t, \mathbf{Q}]$ with $t \in \mathbf{Z}_{\mathbf{h}}^{\times}$(see §8.1). Then $Q\left(\mu ; b_{0}, \mathfrak{b}\right)^{\sigma}=Q\left(\mu ; b_{1}, \mathfrak{b}\right)$ with an element $b_{1} \in F$ such that $\left(t b_{1}-b_{0}\right)_{v} \in \mathfrak{b}_{v}$ for every $v \in \mathbf{h}$.
(2) $D_{\mu}\left(1-\mu ; b_{0}, \mathfrak{b}\right) \in \mathbf{Q}$.

Proof. Take $E_{\mu \mathbf{a}}(z, \mu / 2 ; \lambda)$ as in $\S 18.8$ with $a_{0}=0$ and $\mathfrak{a}=\mathfrak{g}$. Then $\varphi(b) \in$ $\mathbf{Q}_{\mathrm{ab}}$ and for $\sigma$ and $b_{1}$ as above, $\varphi(b)^{\sigma}$ is the quantity obtained from the formula for $\varphi(b)$ with $b_{0}$ replaced by $b_{1}$. Suppose $\mu>1$; exclude the case in which $F=\mathbf{Q}$ and $\mu=2$. Then $\mathcal{D}_{2}=C=0$. Since $\mathcal{D}_{1}(\mu)=N(\mathfrak{b})(-1)^{\mu e} \Gamma(\mu)^{e} Q\left(\mu ; b_{0}, \mathfrak{b}\right)$, Lemma 18.5 shows tht $Q\left(\mu ; b_{0}, \mathfrak{b}\right) \in \mathbf{Q}_{\mathrm{ab}}$ and $Q\left(\mu ; b_{0}, \mathfrak{b}\right)^{\sigma}=Q\left(\mu ; b_{1}, \mathfrak{b}\right)$. If $F=\mathbf{Q}$ and $\mu=2$, then $E_{2 \mathbf{a}}(z, 1 ; \lambda) \in \mathcal{N}_{2}^{1}$, and so we need an analogue of Lemma 18.5 for the elements of $\mathcal{N}_{2}^{1}$. For $f(z)=c(\pi y)^{-1}+\sum_{h \in \mathbf{Q}} a(h) \mathbf{e}(h z) \in \mathcal{N}_{2}^{1}$ and $\tau \in \operatorname{Aut}(\mathbf{C})$ we have $f^{\tau}(z)=c^{\tau}(\pi y)^{-1}+\sum_{h \in \mathbf{Q}} a(h)^{\tau} \mathbf{e}(h z) \in \mathcal{N}_{2}^{1}$ (see $\S 14.15$ ). Now $\mathcal{N}_{2}^{1}$ contains no element, other than 0 , of the form $p(\pi y)^{-1}+q$ with $p, q \in \mathbf{C}$, as can easily be verified. Therefore, modifying the proof of Lemma 18.5, we obtain (1) for $\mu>1$ even when $F=\mathbf{Q}$ and $\mu=2$. Then, from (18.8) we see that $E_{2 \mathbf{a}}(z, 1 ; \lambda)$ belongs to $\pi^{2} \mathcal{N}_{2}^{1}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ in that case. To prove (2) for $\mu>1$, take $E_{\mu \mathbf{a}}(z, \nu / 2 ; \lambda)$ with $b_{0}=0$. Then $\varphi^{\prime}(b) \in \mathbf{Q}$, and hence $D_{\mu \mathbf{a}}\left(\nu-1 ; \kappa_{A}\right) \in \mathbf{Q}$ by Lemma 18.5 (with the modification when $F=\mathbf{Q}$ and $\mu=2$, mentioned above). This proves (2) for $\mu>1$.

Suppose $\mu=1$. Take $a_{0} \in \mathfrak{a}$ and put $\mathfrak{a}^{-1} \mathfrak{d}^{-1}=\mathfrak{c}$; denote by $\kappa^{\prime}$ the characteristic function of $\mathfrak{c}$; define $\kappa_{*}$ as in Lemma 18.2 with $\kappa_{A}$ as $\kappa$. Then $\kappa_{*}=$ $i^{-e} N(\mathfrak{a})^{-1} D_{F}^{-1 / 2} \kappa^{\prime}$. Now the above reasoning applied to the case $\mu=1$ with $a_{0} \in \mathfrak{a}$ shows that

$$
D_{F}^{1 / 2} N(\mathfrak{b})(-2 \pi i)^{-e} D_{\mathbf{a}}\left(1, \kappa_{B}\right)+2^{-e} D_{\mathbf{a}}\left(0, \kappa_{A}\right) \in \mathbf{Q}_{\mathrm{ab}}
$$

In particular, this quantity belongs to $\mathbf{Q}$ if $b_{0} \in \mathfrak{b}$. From the functional equation of Lemma 18.2 we obtain $D_{\mathbf{a}}\left(0, \kappa_{A}\right)=\pi^{-e} D_{\mathbf{a}}\left(1, \kappa_{*}\right)=(\pi i)^{-e} D_{F}^{1 / 2} N(\mathfrak{c}) D_{\mathbf{a}}\left(1, \kappa^{\prime}\right)$. Take $b_{0} \in \mathfrak{b}$ and $\mathfrak{c}=g \mathfrak{b}$ with $g \in F^{\times}$. Then $D_{\mathbf{a}}\left(1, \kappa^{\prime}\right)=g^{-\mathbf{a}} D_{\mathbf{a}}\left(1, \kappa_{B}\right)$, and so

$$
D_{F}^{1 / 2} N(\mathfrak{b})(-2 \pi i)^{-e} D_{\mathbf{a}}\left(1, \kappa_{B}\right)\left\{1+(-g)^{-\mathbf{a}}|g|^{\mathbf{a}}\right\} \in \mathbf{Q} .
$$

Choosing a suitable $g$, we find that $Q\left(1 ; b_{0}, \mathfrak{b}\right) \in \mathbf{Q}$ if $b_{0} \in \mathfrak{b}$. Taking $\mathfrak{c}$ to be $\mathfrak{b}$, we find that $D_{\mathbf{a}}\left(0, \kappa_{A}\right) \in \mathbf{Q}$ if $a_{0} \in \mathfrak{a}$. Then applying $\sigma$ to the Fourier coefficients, we obtain $Q\left(1 ; b_{0}, \mathfrak{b}\right)^{\sigma}=Q\left(1 ; b_{1}, \mathfrak{b}\right)$ for the same reason as in the case $\mu>1$. This proves (1) for $\mu=1$. As for assertion (2) for $\mu=1$, we have seen that $D_{1}\left(0 ; a_{0}, \mathfrak{a}\right) \in \mathbf{Q}$ if $a_{0} \in \mathfrak{a}$. If $a_{0} \notin \mathfrak{a}$, then taking $b_{0}=0$, we immediately see from Lemma 18.5 that $D_{\mathbf{a}}\left(0, \kappa_{A}\right)=2^{e} \mathcal{D}_{2}(0) \in \mathbf{Q}$. This proves (2) for $\mu=1$.

That $Q\left(\mu ; b_{0}, \mathfrak{b}\right) \in \mathbf{Q}_{\mathrm{ab}}$ for $\mu>1$ was proven by Klingen $[\mathrm{K}]$ in a somewhat different formulation.
18.11. Let $\psi$ be a Hecke character of $F$, and $\mathfrak{f}$ the conductor of $\psi$. We understand that the symbol $\psi_{\mathbf{a}}(x) \psi^{*}(x \mathfrak{b})$ for $x=0$ and a $\mathfrak{g}$-ideal $\mathfrak{b}$ denotes $\psi^{*}(\mathfrak{b})$ or 0 according as $\mathfrak{f}=\mathfrak{g}$ or $\mathfrak{f} \neq \mathfrak{g}$. Then, for any fixed $\mathfrak{g}$-ideal $\mathfrak{a}, x \mapsto \psi_{\mathbf{a}}(x) \psi^{*}\left(x \mathfrak{a}^{-1}\right)$ defines a function of $x \in \mathfrak{a} / \mathfrak{a f}$. Indeed, we easily see that the property that $x \mathfrak{a}^{-1}$ is prime to $\mathfrak{f}$ depends only on $x+\mathfrak{a f}$. Take $a \in F_{\mathfrak{h}}^{\times}$so that $\mathfrak{a}=a \mathfrak{g}$. If $x \in \mathfrak{a}$ and $x \mathfrak{a}^{-1}$
is prime to $\mathfrak{f}$, then $\psi_{\mathbf{a}}(x) \psi^{*}\left(x \mathfrak{a}^{-1}\right)=\psi_{\mathbf{a}}(x)\left(\psi_{\mathbf{h}} / \psi_{\mathfrak{f}}\right)\left(x a^{-1}\right)=\psi_{\mathfrak{f}}\left(a x^{-1}\right) \psi\left(a^{-1}\right)$, and clearly this depends only on $x+\mathfrak{a f}$.

We now define the Gauss sum $\mathbf{g}(\psi)$ of $\psi$ by

$$
\mathbf{g}(\psi)= \begin{cases}\psi^{*}(\mathfrak{d}) & \text { if } \mathfrak{f}=\mathfrak{g}  \tag{18.10}\\ \sum_{t \in(\mathfrak{f} \mathfrak{o})^{-1} / \mathfrak{D}^{-1}} \psi_{\mathbf{a}}(t) \psi^{*}(t \mathfrak{f} \mathfrak{d}) \mathbf{e}_{\mathbf{a}}(t) & \text { if } \mathfrak{f} \neq \mathfrak{g}\end{cases}
$$

We note that $|\mathbf{g}(\psi)|^{2}=N(\mathfrak{f})($ see $[\mathrm{S} 97,(\mathrm{~A} 6.3 .2)])$.
18.12. Theorem. Let $0<\mu \in \mathbf{Z}$ and let $\psi$ be a Hecke character of $F$ such that $\psi_{\mathbf{a}}(x)=x^{\mu \mathbf{a}}|x|^{-\mu \mathbf{a}}$. For any $\mathfrak{g}$-ideal $\mathfrak{c}$ put

$$
P_{\mathfrak{c}}(\mu, \psi)=\mathbf{g}(\psi)^{-1}(2 \pi i)^{-\mu e} D_{F}^{1 / 2} L_{\mathfrak{c}}(\mu, \psi)
$$

Then the following assertions hold:
(1) $L(1-\mu, \psi)$ is $P_{\mathfrak{g}}(\mu, \bar{\psi})$ times an element of $\mathbf{Q}^{\times}$.
(2) Both $L_{\mathrm{c}}(1-\mu, \psi)$ and $P_{\mathrm{c}}(\mu, \psi)$ belong to the field generated over $\mathbf{Q}$ by the values of $\psi$, (which is contained in $\mathbf{Q}_{\mathrm{ab}}$, since $\psi$ is of finite order).
(3) For every $\sigma \in \operatorname{Gal}\left(\mathbf{Q}_{\mathrm{ab}} / \mathbf{Q}\right)$ we have $P_{\mathfrak{c}}(\mu, \psi)^{\sigma}=P_{\mathfrak{c}}\left(\mu, \psi^{\sigma}\right)$ and $L_{\mathfrak{c}}(1-$ $\mu, \psi)^{\sigma}=L_{\mathfrak{c}}\left(1-\mu, \psi^{\sigma}\right)$.

Proof. Put $R(s, \psi)=N(\mathfrak{d f})^{s / 2} \Gamma_{(t)}(s) L(s, \psi)$ with $\Gamma_{(t)}$ of Lemma 18.2, where $t=t_{0} \mathbf{a}$ with an integer $t_{0}$ such that $0 \leq t_{0} \leq 1$ and $t_{0}-\mu \in 2 Z$. Then $R(s, \psi)=W_{\psi} R(1-s, \bar{\psi})$ with $W_{\psi}=(-i)^{t_{0} e} N(\mathfrak{f})^{-1 / 2} \mathbf{g}(\psi)$ (see [S97, Theorem A6.2, (A6.3.3)]). Therefore (1) can be verified by a straightforward calculation. Since $L_{\mathfrak{c}}(s, \psi)=L(s, \psi) \prod_{\mathfrak{p} \mid \mathfrak{c}}\left[1-\psi^{*}(\mathfrak{p}) N(\mathfrak{p})^{-s}\right]$, it is sufficient to prove (2) and (3) for $\mathfrak{c}=\mathfrak{g}$. Clearly (2) follows from (3). Thus, in view of (1), we only have to prove that $L(1-\mu, \psi)^{\sigma}=L\left(1-\mu, \psi^{\sigma}\right)$ for every $\sigma \in \operatorname{Gal}\left(\mathbf{Q}_{\mathrm{ab}} / \mathbf{Q}\right)$. Take a complete set of representatives $\left\{\mathfrak{a}_{i}\right\}_{i=1}^{n}$ for the ideal group of $F$ modulo the subgroup $\{p \mathfrak{g} \mid p \in F, N(p)>0\}$. Put $U_{0}=\left\{u \in \mathfrak{g}^{\times} \mid N(u)>0\right\}$ and $U=U_{0} \cap(1+\mathfrak{f})$. Then the ideals $x \mathfrak{a}_{\nu}^{-1}$ for all $\nu$ and all $x \in F^{\times} / U_{0}$ cover the ideal group of $F$ exactly twice, and hence

$$
\begin{aligned}
2\left[U_{0}: U\right] L(s, \psi) & =\sum_{\nu=1}^{m} \sum_{y \in\left(F \times \cap \mathfrak{a}_{\nu}\right) / U} \psi^{*}\left(y \mathfrak{a}_{\nu}^{-1}\right) N\left(y \mathfrak{a}_{\nu}^{-1}\right)^{-s} \\
= & \sum_{\nu=1}^{m} N\left(\mathfrak{a}_{\nu}\right)^{s} \sum_{x \in \mathfrak{a}_{\nu} / \mathfrak{a}_{\nu} \mathfrak{f}} \psi^{*}\left(x \mathfrak{a}_{\nu}^{-1}\right) \sum_{d \in\left(x+\mathfrak{a}_{\nu} \mathfrak{f}\right) / U} \psi^{*}\left(x^{-1} d \mathfrak{g}\right) N(d \mathfrak{g})^{-s}
\end{aligned}
$$

If $d \in x+\mathfrak{a}_{\nu} \mathfrak{f}$ and $\psi^{*}\left(x^{-1} d \mathfrak{g}\right) \neq 0$, then $\psi^{*}\left(x^{-1} d \mathfrak{g}\right)=\psi_{\mathbf{a}}\left(d^{-1} x\right)$, and hence

$$
2\left[U_{0}: U\right] L(s, \psi)=\sum_{\nu=1}^{m} N\left(\mathfrak{a}_{\nu}\right)^{s} \sum_{x \in \mathfrak{a}_{\nu} / \mathfrak{a}_{\nu} \mathfrak{f}} \psi_{\mathbf{a}}(x) \psi^{*}\left(x \mathfrak{a}_{\nu}^{-1}\right) \sum_{d \in\left(x+\mathfrak{a}_{\nu} \mathfrak{f}\right) / U} \psi_{\mathbf{a}}(d) N(d \mathfrak{g})^{-s}
$$

Thus

$$
2 L(1-\mu, \psi)=\left[\mathfrak{g}^{\times}: U_{0}\right] \sum_{\nu=1}^{m} N\left(\mathfrak{a}_{\nu}\right)^{1-\mu} \sum_{x \in \mathfrak{a}_{\nu} / \mathfrak{a}_{\nu} \mathfrak{f}} \psi_{\mathbf{a}}(x) \psi^{*}\left(x \mathfrak{a}_{\nu}^{-1}\right) D_{\mu}\left(1-\mu ; x, \mathfrak{a}_{\nu} \mathfrak{f}\right)
$$

Applying $\sigma$ to this equality, we obtain $L(1-\mu, \psi)^{\sigma}=L\left(1-\mu, \psi^{\sigma}\right)$ as expected.
18.13. For $v \in \mathbf{a}, 0<m \in \mathbf{Z}, k \in \mathbf{Z}^{\mathbf{a}}$, and $0 \leq p \in \mathbf{Z}^{\mathbf{a}}$ we define differential operators $\delta_{m}^{v}$ and $\delta_{k}^{p}$ on $\mathcal{H}$ by

$$
\begin{align*}
& \delta_{m}^{v} f=y^{-m} \frac{\partial}{\partial z_{v}}\left(y_{v}^{m} f\right)=\left(\frac{m}{2 i y_{v}}+\frac{\partial}{\partial z_{v}}\right) f,  \tag{18.11}\\
& \delta_{k}^{p}=\prod_{v \in \mathbf{a}} \delta_{k_{v}+2 p_{v}-2}^{v} \cdots \delta_{k_{v}+2}^{v} \delta_{k_{v}}^{v} \tag{18.12}
\end{align*}
$$

where $f$ is a function on $\mathcal{H}$. As explained in $\S 12.17, \delta_{k}^{p}$ is a special case of $D_{\rho}^{p}$, and maps $\mathcal{N}_{k}^{q}$ into $\mathcal{N}_{k+2 p}^{q+p}$ for $0 \leq q \in \mathbf{Z}^{\mathbf{a}}$. Moreover, $(\pi i)^{-|p|} \delta_{k}^{p}$ maps $\mathcal{N}_{k}^{q}(W)$ into $\mathcal{N}_{k+2 p}^{q+p}(W)$ for every subfield $W$ of $\mathbf{C}$ containing the Galois closure of $F$ over $\mathbf{Q}$. This is a special case of Theorem 14.12 (4), and in fact, can be verified by a direct calculation employing the expression (14.21) for the elements of $\mathcal{N}_{k}^{q}$. Now we have

$$
\begin{gather*}
\delta_{k}^{p}\left(f \|_{k} \alpha\right)=\left(\delta_{k}^{p} f\right) \|_{k+2 p} \alpha \quad \text { for every } \quad \alpha \in G  \tag{18.13}\\
\delta_{k}^{p}\left(y^{r}\right)=(2 i)^{-|p|} y^{r-p} \prod_{v \in \mathbf{a}} \frac{\Gamma\left(k_{v}+r_{v}+p_{v}\right)}{\Gamma\left(k_{v}+r_{v}\right)} \quad\left(r \in \mathbf{C}^{\mathbf{a}}\right) \tag{18.14}
\end{gather*}
$$

The former is a special case of Proposition 12.10 (2) as noted in $\S 12.17$; the latter can be verified by a direct calculation. Next, for $h$ in the set $H$ of $\S 18.3$ we have

$$
\begin{align*}
& \delta_{k}^{p}\left\{j_{h}(z)^{-k}\left|j_{h}(z)\right|^{k-2 s \mathbf{a}} y^{s \mathbf{a}-k / 2}\right\}  \tag{18.15}\\
= & (2 i)^{-|p|} y^{s \mathbf{a}-\ell / 2} j_{h}(z)^{-\ell}\left|j_{h}(z)\right|^{\ell-2 s \mathbf{a}} \prod_{v \in \mathbf{a}} \frac{\Gamma\left(s+k_{v} / 2+p_{v}\right)}{\Gamma\left(s+k_{v} / 2\right)}
\end{align*}
$$

with $\ell=k+2 p$. Indeed, we can find $\alpha \in G$ such that $j_{\alpha}=j_{h}$. Then the left-hand side of (18.15) is $\delta_{k}^{p}\left(y^{s \mathbf{a}-k / 2} \|_{k} \alpha\right.$ ), which equals $\left(\delta_{k}^{p} y^{s \mathbf{a}-k / 2}\right) \|_{k+2 p} \alpha$ by (18.13). From this and (18.14) we obtain (18.15). Thus, if $\operatorname{Re}(s)$ is sufficiently large,

$$
\begin{equation*}
\delta_{k}^{p} E_{k}(z, s ; \lambda)=(2 i)^{-|p|} E_{k+2 p}(z, s ; \lambda) \prod_{v \in \mathbf{a}} \frac{\Gamma\left(s+k_{v} / 2+p_{v}\right)}{\Gamma\left(s+k_{v} / 2\right)} \tag{18.16}
\end{equation*}
$$

By analytic continuation, the equality holds for every $s$.
18.14. Theorem. Let $\Phi$ be the Galois closure of $F$ over $\mathbf{Q}$, and $\lambda$ a $\mathbf{Q}_{\mathrm{ab}}$ valued element of $\mathcal{S}\left(\left(F_{2}^{1}\right)_{\mathbf{h}}\right)$; let $k$ be an element of $\mathbf{Z}^{\mathbf{a}}$ such that $k_{v} \geq 1$ for every $v \in \mathbf{a}$ and $k_{v}-k_{v^{\prime}} \in 2 \mathbf{Z}$ for every $v, v^{\prime} \in \mathbf{a}$. Then, for every $\mu \in \mathbf{Z}$ such that $\mu-k_{v} \in 2 \mathbf{Z}$ and $2-k_{v} \leq \mu \leq k_{v}$ for every $v \in \mathbf{a}$, the function $E_{k}(z, \mu / 2 ; \lambda)$ belongs to $\pi^{\alpha} \mathcal{N}_{k}^{t}\left(\Phi \mathbf{Q}_{\mathrm{ab}}\right)$, where $\alpha=(1 / 2) \sum_{v \in \mathbf{a}}\left(k_{v}+\mu\right)$, and $t=(k-|\mu-1| \mathbf{a}-\mathbf{a}) / 2$, except when $F=\mathbf{Q}$ and $|\mu-1|=1$, in which case $t=k / 2$.

Proof. First suppose that $\mu>0$; put $p=(k-\mu \mathbf{a}) / 2$. Then $0 \leq p \in \mathbf{Z}^{\mathbf{a}}$, and (18.16) shows that

$$
\delta_{\mu \mathbf{a}}^{p} E_{\mu \mathbf{a}}(z, \mu / 2 ; \lambda)=(2 i)^{-|p|} c E_{k}(z, \mu / 2 ; \lambda)
$$

with $c \in \mathbf{Q}^{\times}$. Therefore our assertion follows from Theorem 18.9 (1) combined with the fact that $\pi^{-|p|} \delta_{\mu \mathbf{a}}^{p}$ sends $\mathcal{N}_{\mu \mathbf{a}}^{q}\left(\Phi \mathbf{Q}_{\mathrm{ab}}\right)$ into $\mathcal{N}_{k}^{q+p}\left(\Phi \mathbf{Q}_{\mathrm{ab}}\right)$. Next, take an integer $\nu$ such that $2-k_{v} \leq \nu \leq 0$ and $\nu-k_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$. Put $\mu=2-\nu$ and $p=(k-\mu \mathbf{a}) / 2$. Then $0 \leq p \in \mathbf{Z}^{\mathbf{a}}$, and (18.16) shows that

$$
\delta_{\mu \mathbf{a}}^{p} E_{\mu \mathbf{a}}(z, \nu / 2 ; \lambda)=(2 i)^{-|p|} c^{\prime} E_{k}(z, \nu / 2 ; \lambda)
$$

with $c^{\prime} \in \mathbf{Q}^{\times}$. By virtue of Theorem 18.9 (2), we obtain the desired conclusion for the same reason as in the case $\mu>0$.
18.15. We now take a totally imaginary quadratic extension $K$ of $F$, and denote by $\mathfrak{r}$ the maximal order of $K$. We consider the $L$-function $L(s, \chi)$ of a Hecke character $\chi$ of $K$ such that $\chi_{\mathbf{a}}(x)=\prod_{\sigma \in J_{K}}\left(x^{\sigma} /\left|x^{\sigma}\right|\right)^{\mu_{\sigma}}$ with $\mu \in I_{K}$, where $J_{K}$
and $I_{K}$ are as in §11.3. Then we can easily find a CM-type $\tau=\sum_{v \in \mathbf{a}} \tau_{v}$ of $K$ such that

$$
\begin{equation*}
\chi_{\mathbf{a}}(x)=\prod_{v \in \mathbf{a}}\left(x^{\tau_{v}} /\left|x^{\tau_{v}}\right|\right)^{m_{v}} \tag{18.17}
\end{equation*}
$$

with $0 \leq m_{v} \in \mathbf{Z}$. Fixing a CM-type $(K, \tau)$ and $m \in \mathbf{Z}^{\mathbf{a}}, \geq 0$, put

$$
\begin{equation*}
L_{m}(s, \ell)=\left[\mathfrak{r}^{\times}: U\right]^{-1} \sum_{\alpha \in K^{\times} / U} \ell(\alpha) \alpha^{-m}|\alpha|^{m-2 s \mathbf{a}} \tag{18.18}
\end{equation*}
$$

Here $\ell \in \mathcal{S}\left(K_{\mathbf{h}}\right), \alpha^{-m}|\alpha|^{m-2 s \mathbf{a}}=\prod_{v \in \mathbf{a}}\left(\alpha^{\tau_{v}}\right)^{-m_{v}}\left|\alpha^{\tau_{v}}\right|^{m_{v}-2 s}, U$ is a subgroup of $\mathfrak{r}^{\times}$of finite index such that $\ell(u \alpha)=\ell(\alpha)$ and $u^{-m}|u|^{m}=1$ for every $u \in U$. Clearly such a $U$ exists, and $L_{m}(s, \ell)$ does not depend on the choice of $U$; we easily see that the sum is convergent for $\operatorname{Re}(s)>1$.
18.16. Theorem. The notation being as above, suppose $m \neq 0$; then the following assertions hold:
(1) $L_{m}(s, \ell) \prod_{v \in \mathbf{a}} \Gamma\left(s+m_{v} / 2\right)$ can be continued as an entire function of $s$ to the whole $\mathbf{C}$.
(2) Suppose that $\ell$ is $\overline{\mathbf{Q}}$-valued and $m_{v}-m_{v^{\prime}} \in 2 \mathbf{Z}$ for every $v, v^{\prime} \in \mathbf{a}$; let $\mu$ be an integer such that $2-m_{v} \leq \mu \leq m_{v}$ and $\mu-m_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$. Then $L_{m}(\mu / 2, \ell) \in \pi^{\alpha} p_{K}\left(\sum_{v \in \mathbf{a}} m_{v} \tau_{v}, \tau\right) \overline{\overline{\mathbf{Q}}}$, where $\alpha=(1 / 2) \sum_{v \in \mathbf{a}}\left(m_{v}+\mu\right)$ and $p_{K}$ is the symbol of §11.3.
(3) If $\chi$ is a Hecke character of type (18.17), then $L(\mu / 2, \chi) \in \pi^{\alpha} p_{K}\left(\sum_{v \in \mathbf{a}} m_{v} \tau_{v}\right.$, $\tau) \overline{\mathbf{Q}}$ for $\mu$ and $\alpha$ as in (2).

Proof. Assertion (1) will be proven in §A7.4. To prove (2), we first assume that $\mu \geq 1$. Take $w_{0} \in K$ so that $\operatorname{Im}\left(w_{0}^{\tau_{v}}\right)>0$ for every $v \in \mathbf{a}$, and put $w=$ $\left(w_{0}^{\tau_{v}}\right)_{v \in \mathbf{a}}$. By Proposition 4.14, $w$ is a CM-point of $\mathcal{H}$. Define an $F$-linear bijection $g: F_{2}^{1} \rightarrow K$ by $g(c, d)=c w_{0}+d$, and put $\lambda=\ell \circ g$. Then $\lambda \in \mathcal{S}\left(\left(F_{2}^{1}\right)_{\mathbf{h}}\right)$ and $j_{h}(w)^{m}=g(h)^{m}$ for $h \in F_{2}^{1}, \neq 0$, and hence

$$
\begin{aligned}
L_{m}(s, \ell) & =\left[\mathfrak{r}^{\times}: U\right]^{-1} \sum_{h \in H / U} \lambda(h) j_{h}(w)^{-m}\left|j_{h}(w)\right|^{m-2 s \mathbf{a}} \\
& =\left[\mathfrak{r}^{\times}: \mathfrak{g}^{\times}\right]^{-1} y^{m / 2-s \mathbf{a}} E_{m}(w, s ; \lambda)
\end{aligned}
$$

where $y=\operatorname{Im}(w)$. (Notice that we may take $U \subset \mathfrak{g}^{\times}$in (18.18).) Given $\mu$ and $\alpha$ as in (2), we have $E_{m}(z, \mu / 2 ; \lambda) \in \pi^{\alpha} \mathcal{N}_{m}^{t}(\overline{\mathbf{Q}})$ by Theorem 18.14, and hence its value at $w$ belongs to $\pi^{\alpha} \mathfrak{P}_{m}(w) \overline{\mathbf{Q}}$. Since $y$ is algebraic and $\mathfrak{P}_{m}(w)=p_{K}\left(\sum_{v \in \mathbf{a}} m_{v} \tau_{v}, \tau\right)$ by Proposition 11.18, we obtain (2). To prove (3), let $\mathfrak{h}$ be the conductor of $\chi$, and $A$ a complete set of representatives for the ideal classes of $K$ modulo $\mathfrak{h}$, consisting of integral ideals prime to $\mathfrak{h}$. Then

$$
L(s, \chi)=\sum_{\mathfrak{a} \in A} \chi^{*}(\mathfrak{a})^{-1} N(\mathfrak{a})^{s} \sum_{\alpha \in W_{\mathfrak{a}} / U_{\mathfrak{h}}} \chi^{*}(\alpha \mathfrak{r}) N(\alpha \mathfrak{r})^{-s}
$$

with $W_{\mathfrak{a}}=\mathfrak{a} \cap(1+\mathfrak{h}) \cap K^{\times}, 1+\mathfrak{h}=\{1+x \mid x \in \mathfrak{h}\}$, and $U_{\mathfrak{h}}=\mathfrak{r}^{\times} \cap(1+\mathfrak{h})$. Let $\ell_{\mathfrak{a}}$ be the characteristic function of $\mathfrak{a} \cap(1+\mathfrak{h})$. Since $\chi^{*}(\alpha \mathfrak{r})=\chi_{\mathbf{a}}(\alpha)^{-1}=\alpha^{-m}|\alpha|^{m}$ for $\alpha \in 1+\mathfrak{h}$, we have $u^{-m}|u|^{m}=1$ for $u \in U_{\mathfrak{h}}$, and hence

$$
L(s, \chi)=\left[\mathfrak{r}^{\times}: U_{\mathfrak{h}}\right] \sum_{\mathfrak{a} \in A} \chi^{*}(\mathfrak{a})^{-1} N(\mathfrak{a})^{s} L_{m}\left(s, \ell_{\mathfrak{a}}\right) .
$$

Now $\chi^{*}(\mathfrak{a}) \in \overline{\mathbf{Q}}$ by Lemma 17.11. Therefore (3) follows from (2).

## CHAPTER V

## ZETA FUNCTIONS ASSOCIATED WITH HECKE EIGENFORMS

## 19. Formal Euler products and generalized Möbius functions

19.1. Our setting is the same as in Section 16; thus we consider only Cases SP and UT. We do not consider $G_{1}=S U\left(\eta_{n}\right)$ with $K \neq F$ in this section. Let us first make the following notational convention: For $v \in \mathbf{h}$ and a subgroup $X$ of $G_{\mathbf{A}}$ (resp. $G L_{n}(K)_{\mathbf{A}}$ ) we put $X_{v}=X \cap G_{v}$ (resp. $X_{v}=X \cap G L_{n}\left(K_{v}\right)$ ). In fact, whenever this notation is used, $X_{v}$ is the projection of $X$ to $G_{v}$. We now take a fractional ideal $\mathfrak{b}$ and an integral ideal $\mathfrak{c}$ in $F$, and consider the subgroup $D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]$ of $G_{\mathbf{A}}$ defined in $\S 16$. We also take a divisor $\mathfrak{e}$ of $\mathfrak{c}$ such that $\mathfrak{e}^{-1} \mathfrak{c}+\mathfrak{e}=\mathfrak{g}$, and throughout this section we denote by $C$ the subgroup of $G_{\mathbf{A}}$ defined by

$$
\begin{equation*}
C=\left\{x \in D\left[\mathfrak{b}^{-1} \mathfrak{e}, \mathfrak{b} \mathfrak{c}\right] \mid a_{x}-1 \prec \mathfrak{r e}\right\} . \tag{19.1}
\end{equation*}
$$

Recall the equality $G_{\mathbf{A}}=P_{\mathbf{A}} D\left[\mathfrak{b}^{-1}, \mathfrak{b}\right]$ as stated in (16.10). It may also be emphasized that $\mathfrak{b}, \mathfrak{c}$, and $\mathfrak{e}$ are in $F$.

To define Hecke operators, put

$$
\begin{gather*}
E=\prod_{v \in \mathbf{h}} G L_{n}\left(\mathfrak{r}_{v}\right), \quad B=\left\{x \in G L_{n}(K)_{\mathbf{h}} \mid x \prec \mathfrak{r}\right\},  \tag{19.2a}\\
B^{\prime}=\{x \in B \mid x-1 \prec \mathfrak{r e}\}, \quad E^{\prime}=B^{\prime} \cap E,  \tag{19.2b}\\
\mathfrak{X}=C Q(\mathfrak{e}) C, \quad Q(\mathfrak{e})=\left\{\operatorname{diag}[\widehat{r}, r] \mid r \in B^{\prime}\right\} . \tag{19.2c}
\end{gather*}
$$

By [S97, Proposition 5.10 or 7.8 ] (see also Remark 16.12 (III)) we have

$$
\begin{equation*}
G_{\mathbf{A}}=D\left[\mathfrak{b}^{-1}, \mathfrak{b}\right] Q(\mathfrak{g}) D\left[\mathfrak{b}^{-1}, \mathfrak{b}\right] . \tag{19.3}
\end{equation*}
$$

Thus $G_{v}=C_{v} Q(\mathfrak{e})_{v} C_{v}$ for every $v \nmid \mathfrak{c}$, and so $\mathfrak{X}_{v}=G_{v}$ for such a $v$; clearly $Q(\mathfrak{e})_{v} \subset C_{v}$ and $\mathfrak{X}_{v}=C_{v}$ for $v \mid \mathfrak{e}$. We are primarily interested in the cases $\mathfrak{e}=\mathfrak{g}$ and $\mathfrak{e}=\mathfrak{c}$; if $\mathfrak{e}=\mathfrak{c}$, the group is essentially a principal congruence subgroup; if $\mathfrak{e}=\mathfrak{g}$, it is an analogue of the group $\Gamma_{0}(N)$ in the elliptic modular case.
19.2. Lemma. (1) $\mathfrak{X} \subset C P_{\mathbf{h}}$.
(2) Let $J_{v}=C_{v} \cap P_{v}$. If $v \mid \mathfrak{c}$ and $\sigma \in Q(\mathfrak{e})_{v}$, then $C_{v} \sigma C_{v} \cap P_{v}=J_{v} \sigma J_{v}$ and $C_{v} \sigma C_{v}=C_{v} \sigma J_{v}$.
(3) Let $\sigma=\operatorname{diag}[\widehat{q}, q] \in Q(\mathfrak{e})_{v}$ with $v \mid \mathbf{c}, v \nmid e$. Then $C_{v} \sigma C_{v}=\bigsqcup_{d, b} C_{v}\left[\begin{array}{cc}\widehat{d} & \widehat{d b} \\ 0 & d\end{array}\right]$ with $d \in E_{v} \backslash E_{v} q E_{v}$ and $b \in S\left(\mathfrak{b}^{-1}\right)_{v} / d^{*} S\left(\mathfrak{b}^{-1}\right)_{v} d$, where $S()$ is as in (16.1b).

Proof. To prove (1), it is sufficient to show that $\sigma C_{v} \subset C_{v} P_{v}$ for every $\sigma=$ $\operatorname{diag}[\widehat{q}, q] \in Q(\mathfrak{e})_{v}, v \in \mathbf{h}$. Since $C_{v} P_{v}=G_{v}$ if $v \nmid c$ and $\mathfrak{X}_{v}=C_{v}$ if $v \mid \mathfrak{e}$, we may assume that $v \mid \mathfrak{c}$ and $v \nmid \mathfrak{e}$. Put $A_{v}=D\left[\mathfrak{b}^{-1}, \mathfrak{b}\right]_{v}$. Let $\alpha \in C_{v}$. Since $G_{v}=A_{v} P_{v}$,
we have $\sigma \alpha=\beta^{-1} \pi$ with $\beta \in A_{v}$ and $\pi \in P_{v}$. Since $v \mid \mathfrak{c}$, both $a_{\alpha}$ and $d_{\alpha}$ belong to $G L_{n}\left(\mathfrak{r}_{v}\right)$. Also $0=c_{\pi}=c_{\beta} \widehat{q} a_{\alpha}+d_{\beta} q c_{\alpha}$, we have $c_{\beta}=-d_{\beta} q c_{\alpha} a_{\alpha}^{-1} q^{*} \prec \mathfrak{r}_{v} \mathfrak{b}_{v} \mathfrak{c}_{v}$, and hence $\beta \in C_{v}$, which proves (1). To prove (2), we may again assume that $v \nmid e$, since $C_{v} \sigma C_{v}=C_{v}$ and $\sigma \in J_{v}$ if $v \mid \mathfrak{e}$. Given $\alpha=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right] \in C_{v}, v \mid \mathfrak{c}$, put $\beta=\left[\begin{array}{cc}d^{*} & -b^{*} \\ 0 & d^{-1}\end{array}\right]$. Then $\beta \in J_{v}$ and $\beta \alpha=\left[\begin{array}{ll}1 & 0 \\ s & 1\end{array}\right]$ with $s \in S_{v}$. Now let $\sigma=$ $\operatorname{diag}[\widehat{q}, q] \in Q(\mathfrak{e})_{v}$ and $\pi=\alpha_{1} \sigma \alpha_{2}^{-1} \in P_{v}$ with $\alpha_{i} \in C_{v}$. Applying the result just proved to $\alpha_{i}$, we obtain elements $\beta_{i} \in J_{v}$ such that $\beta_{i} \alpha_{i}=\left[\begin{array}{cc}1 & 0 \\ s_{i} & 1\end{array}\right]$ with $s_{i} \in S_{v}$. Then $\beta_{1} \pi \beta_{2}^{-1}=\left[\begin{array}{ll}\widehat{q} & 0 \\ e & q\end{array}\right]$ with some $e \in\left(K_{v}\right)_{n}^{n}$. Since $\beta_{1} \pi \beta_{2}^{-1} \in P_{v}$, we see that $e=0$, which shows that $\pi \in J_{v} \sigma J_{v}$, and hence $C_{v} \sigma C_{v} \cap P_{v}=J_{v} \sigma J_{v}$. Next, take any $\gamma, \varepsilon \in C_{v}$. By (1), $\gamma \sigma \varepsilon=\xi \zeta$ with $\xi \in C_{v}$ and $\zeta \in P_{v}$. Then $\zeta \in C_{v} \sigma C_{v} \cap P_{v}=J_{v} \sigma J_{v}$, and hence $\gamma \sigma \varepsilon \in \xi J_{v} \sigma J_{v} \subset C_{v} \sigma J_{v}$, from which we obtain $C_{v} \sigma C_{v}=C_{v} \sigma J_{v}$. As for (3), since $C_{v} \sigma C_{v}=C_{v} \sigma J_{v}$, each coset of $C_{v} \backslash C_{v} \sigma C_{v}$ has a representative in $\sigma J_{v}$. Then our assertion can be verified in a straightforward way.
19.3. Lemma. The set $\mathfrak{X}$ is closed under multiplication.

Proof. It is sufficient to prove that $\mathfrak{X}_{v}$ is closed under multiplication for every $v \in \mathbf{h}$; again we may assume that $v \mid \mathbf{c}$ and $v \nmid e$. Let $C_{v} \sigma C_{v}=\bigsqcup_{d, b} C_{v}[[d, b]]$ with $\sigma \in Q(\mathfrak{e})_{v}$ as in Lemma $19.2(3)$, where $[[d, b]]=\left[\begin{array}{cc}\widehat{d} & \widehat{d b} \\ 0 & d\end{array}\right]$; similarly let $C_{v} \tau C_{v}=$ $\bigsqcup_{e, g} C_{v}[[e, g]]$ for $\tau \in Q(\mathfrak{e})_{v}$. Then $C_{v} \sigma C_{v} \tau C_{v}=\bigsqcup_{d, b, e, g} C_{v}[[d, b]][[e, g]]=\bigsqcup_{d, b, e, g}$ $C_{v} \operatorname{diag}[\widehat{d} \widehat{e}, d e][[1, x]]$, where $x=g+e^{*} b e$. Since $[[1, x]] \in C_{v}$ and $\operatorname{diag}[\widehat{d} \widehat{e}, d e] \in$ $Q(e){ }_{v}$, we obtain the desired result.
19.4. Lemma. Employing the symbols of (16.1a, b, c), put $L_{0}=\mathfrak{r}_{1}^{n}$ and

$$
\begin{aligned}
W=\{(g, h) & \left.\in B^{\prime} \times B^{\prime} \mid g L_{0}+h L_{0}=L_{0}, h_{v} \in E_{v}^{\prime} \text { for every } v \mid \mathfrak{c}\right\} \\
S^{\prime} & =\left\{\sigma \in S_{\mathbf{h}} \mid \sigma_{v} \in S\left(\mathfrak{b}^{-1} \mathfrak{e}\right)_{v} \text { for every } v \mid \mathfrak{c}\right\}
\end{aligned}
$$

Then a complete set of representatives for $C \backslash \mathfrak{X}$ can be given by

$$
\left\{\left.\left[\begin{array}{cc}
g^{-1} h & g^{-1} \sigma \widehat{h} \\
0 & g^{*} \widehat{h}
\end{array}\right] \right\rvert\,(g, h) \in E^{\prime} \backslash W /\left(E^{\prime} \times 1\right), \quad \sigma \in S^{\prime} / g S_{\mathbf{h}}\left(\mathfrak{b}^{-1} \mathfrak{e}\right) g^{*}\right\}
$$

where we let $e \in E^{\prime}$ and $(f, 1) \in E^{\prime} \times 1$ act on $W$ by $e(g, h)(f, 1)=(e g f$, eh $)$.
Proof. This is essentially the local problem of finding $C_{v} \backslash \mathfrak{X}_{v}$. This is trivial if $v \mid \mathbf{e}$. If $v \nmid \mathfrak{c}$, then the question is about $C_{v} \backslash G_{v}$. Since $G_{v}=C_{v} P_{v}$, we can take representatives from $\left(C_{v} \cap P_{v}\right) \backslash P_{v}$. Now $\left(C_{v} \cap P_{v}\right) \backslash P_{v}$ can be given by $\left[\begin{array}{cc}a & s \widehat{a} \\ 0 & \widehat{a}\end{array}\right]$ with $a \in E_{v} \backslash G L_{n}\left(K_{v}\right)$ and $s \in S_{v} / S\left(\mathfrak{b}^{-1}\right)_{v}$, as noted in [S97, p.131, last 4 lines]. Also in [S97, Lemma 16.4] we showed that the map $(g, h) \mapsto g^{-1} h$ gives a bijection of $E_{v} \backslash W_{v}$ onto $G L_{r}\left(K_{v}\right)$. Therefore the desired fact can easily be verified. If $v \mid c$ and $v \nmid e$, then by Lemma $19.2(3), C_{v} \backslash \mathfrak{X}_{v}$ can be given by $\left[\begin{array}{cc}\widehat{d} & \widehat{d b} \\ 0 & d\end{array}\right]$ with $d \in E_{v} \backslash B_{v}$ and $b \in S\left(\mathfrak{b}^{-1}\right)_{v} / d^{*} S\left(\mathfrak{b}^{-1}\right)_{v} d$. This time $W_{v}=B_{v} \times E_{v}$ and $E_{v} \backslash W_{v} /\left(E_{v} \times 1\right)=$ $\left(B_{v} / E_{v}\right) \times 1$, and so the same conclusion holds in a much simpler way.
19.5. We shall state some of our results in terms of formal Dirichlet series of the form $\sum_{\mathfrak{a} \subset \mathfrak{r}} c(\mathfrak{a})[\mathfrak{a}]$. Here $c(\mathfrak{a}) \in \mathbf{C}$ and $\{[\mathfrak{a}]\}$ is a system of formal multiplicative symbols defined for the fractional ideals $\mathfrak{a}$ in $K$ as follows: the [ $\mathfrak{p ]}$ for the prime ideals $\mathfrak{p}$ are independent indeterminates; $[\mathfrak{r}]=1$ and $[\mathfrak{a b}]=[\mathfrak{a}][\mathfrak{b}]$. Substituting $\varphi(\mathfrak{a}) N(\mathfrak{a})^{-s}$ with an ideal-character $\varphi$ for [a], we obtain an ordinary Dirichlet series. Hereafter we denote by $\sum_{\mathfrak{a}}$ the sum over all the integral ideals $\mathfrak{a}$ in $K$.

Given $\xi \in G_{\mathbf{A}}$, take $q \in B$ so that $\xi \in D\left[\mathfrak{b}^{-1}, \mathfrak{b}\right] \operatorname{diag}\left[q^{-1}, q^{*}\right] D\left[\mathfrak{b}^{-1}, \mathfrak{b}\right]$, which is feasible by virtue of (19.3). Then we put

$$
\begin{equation*}
\nu_{\mathfrak{b}}(\xi)=\operatorname{det}(q) \mathfrak{r} \tag{19.4}
\end{equation*}
$$

If $\beta=\operatorname{diag}\left[1_{n}, b_{0}^{-1} 1_{n}\right]$ with an lement $b_{0} \in F_{\mathrm{h}}^{\times}$such that $\mathfrak{b}=b_{0} \mathfrak{g}$, then

$$
\begin{equation*}
\nu_{\mathfrak{b}}(\xi)=\nu_{0}\left(\beta \xi \beta^{-1}\right) \tag{19.5}
\end{equation*}
$$

with $\nu_{0}$ defined in $\S 1.10$. Let $\mathfrak{d}$, $\delta$, and $\widetilde{S}$ be as in $\S 16.1$. Given $\zeta \in \widetilde{S}$, we define a formal Dirichlet series $\alpha_{\mathrm{c}}^{0}(\zeta)$ by

$$
\begin{equation*}
\alpha_{c}^{0}(\zeta)=\prod_{v \nmid c} \alpha_{v}^{0}(\zeta), \quad \alpha_{v}^{0}(\zeta)=\sum_{\tau \in S_{v} / S(\mathrm{r}) v} \mathbf{e}_{v}^{n}\left(-\delta_{v}^{-1} \zeta \tau\right)\left[\nu_{0}(\tau)\right] \tag{19.6}
\end{equation*}
$$

This can be obtained by substituting $\left(\tau,\left[\nu_{0}(\tau)\right]\right)$ for $\left(\sigma, \nu[\sigma]^{-s}\right)$ in (16.7a). Clearly

$$
\begin{equation*}
\alpha_{\mathrm{c}}^{0}\left(c \gamma \zeta \gamma^{*}\right)=\alpha_{\mathrm{c}}^{0}(\zeta) \text { if } c \in \prod_{v \in \mathrm{~h}} \mathfrak{g}_{v}^{\times} \text {and } \gamma \in E \tag{19.7}
\end{equation*}
$$

19.6. Lemma. Let $S^{\prime}$ be as in Lemmas 19.4, and $b_{0}$ be as above. Let $\zeta \in S_{\mathbf{h}}$ and $g \in B$; suppose $g^{*} \zeta g \in \mathfrak{b} \mathfrak{e}^{-1} \widetilde{S}$. Then

$$
\sum_{\tau \in X} \mathbf{e}_{\mathbf{h}}^{n}\left(-\delta^{-1} \zeta \tau\right)\left[\nu_{0}\left(b_{0} \tau\right)\right] \neq 0 \quad \text { for } \quad X=S^{\prime} / g S_{\mathbf{h}}\left(\mathfrak{b}^{-1} \mathfrak{e}\right) g^{*}
$$

only if $\zeta \in \mathfrak{b e}^{-1} \widetilde{S}$, in which case the sum equals $|\operatorname{det}(g)|_{K}^{-\kappa} \alpha_{\mathfrak{c}}^{0}\left(b_{0}^{-1} \zeta\right)$, where $|x|_{K}$ denotes the idele norm of $x \in K_{\mathbf{A}}^{\times}$, and $\kappa=n+1$ in Case $S P$ and $\kappa=n$ in Case UT.

Proof. Recall that $\nu_{0}(\tau)$ depends only on $\tau \bmod S_{\mathrm{h}}(\mathfrak{r})$. Change $\tau$ for $\tau+$ $\gamma$ with $\gamma \in S_{\mathbf{h}}\left(\mathfrak{b}^{-1} \mathfrak{e}\right)$. Then the sum in question is multiplied by the factor $\mathbf{e}_{\mathrm{h}}^{n}\left(-\delta^{-1} \zeta \gamma\right)$, which is nonzero for some such $\gamma$ if $\zeta \notin \mathfrak{b e} \mathfrak{e}^{-1} \widetilde{S}$. Thus the sum is nonzero only if $\zeta \in \mathfrak{b e}{ }^{-1} \widetilde{S}$, in which case the sum is $\alpha_{\mathfrak{c}}^{0}\left(b_{0}^{-1} \zeta\right)$ times $\left[S_{\mathbf{h}}\left(\mathfrak{b}^{-1} \mathfrak{e}\right)\right.$ : $\left.g S_{\mathbf{h}}\left(\mathfrak{b}^{-1} \mathfrak{e}\right) g^{*}\right]$. The last number is $|\operatorname{det}(g)|_{K}^{-\kappa}$ as noted in [S97, Lemma 13.2]. This proves our lemma.
19.7. We consider the $\mathbf{Q}$-algebra $\mathfrak{R}(C, \mathfrak{X})$ spanned by the $C \xi C$ for all $\xi \in \mathfrak{X}$ over $\mathbf{Q}$, with the law of multiplication defined as usual (see [S97, Section 11]). This is meaningful because of Lemma 19.3. Similarly we can consider $\mathfrak{R}\left(C_{v}, \mathfrak{X}_{v}\right)$ for each $v \in \mathbf{h}$. These algebras are commutative. Indeed, the commutativity of $\mathfrak{R}(C, \mathfrak{X})$ can be reduced to that of $\mathfrak{R}\left(C_{v}, \mathfrak{X}_{v}\right)$ as can easily be seen. If $v \mid \boldsymbol{e}$, then $\mathfrak{R}\left(C_{v}, \mathfrak{X}_{v}\right)$ is just $\mathbf{Q} \cdot C_{v} 1 C_{v}$. The commutativity of $\mathfrak{R}\left(C_{v}, \mathfrak{X}_{v}\right)$ for $v \nmid e$ follows from the existence of an injection $\omega_{v}$ into a commutative ring as will be shown in Theorem 19.8 below.

We now define formal Dirichlet series $\mathfrak{T}$ and $\mathfrak{T}_{v}$ with coefficients in $\mathfrak{R}(C, \mathfrak{X})$ and $\mathfrak{R}\left(C_{v}, \mathfrak{X}_{v}\right)$ by

$$
\begin{equation*}
\mathfrak{T}=\sum_{\xi \in C \backslash \mathfrak{X} / C} C \xi C\left[\nu_{\mathfrak{b}}(\xi)\right], \quad \mathfrak{T}_{v}=\sum_{\xi \in C_{v} \backslash \mathfrak{X}_{v} / C_{v}} C_{v} \xi C_{v}\left[\nu_{\mathfrak{b}}(\xi)\right] \tag{19.8}
\end{equation*}
$$

If $\xi=\operatorname{diag}\left[r^{-1}, r^{*}\right]$ with $r \in B^{\prime}$, then looking at the elementary divisors of $\xi$ and $r$, we easily see that $C \xi C$ determines $E^{\prime} r E^{\prime}$ and vice versa. Therefore we have

$$
\begin{equation*}
\mathfrak{T}=\sum_{r \in E^{\prime} \backslash B^{\prime} / E^{\prime}} C \operatorname{diag}\left[r^{-1}, r^{*}\right] C[\operatorname{det}(r) \mathfrak{r}] \tag{19.9}
\end{equation*}
$$

For an integral $\mathfrak{r}$-ideal $\mathfrak{a}$ we denote by $T(\mathfrak{a})$ the sum of all the different $C \xi C$ with $\xi \in \mathfrak{X}$ such that $\nu_{\mathfrak{b}}(\xi)=\mathfrak{a}$. Then clearly

$$
\begin{equation*}
\mathfrak{T}=\prod_{v \in \mathbf{h}} \mathfrak{T}_{v}=\sum_{\mathfrak{a}} T(\mathfrak{a})[\mathfrak{a}] . \tag{19.10}
\end{equation*}
$$

We have $\mathfrak{T}_{v}=1$ if $v \mid \mathfrak{e}$, since $C_{v}=\mathfrak{X}_{v}$ for such a $v$.
19.8. Theorem. Let $t_{1}, \ldots, t_{m}$ be $m$ indeterminates, where $m=2 n$ if $K \neq$ $F$ and $v$ splits in $K$, and $m=n$ otherwise. Then for each $v \in \mathbf{h}$ prime to $\mathfrak{e}$ there exists a Q-linear ring-injection

$$
\begin{equation*}
\omega: \mathfrak{R}\left(C_{v}, \mathfrak{X}_{v}\right) \rightarrow \mathbf{Q}\left[t_{1}, \ldots, t_{m}, t_{1}^{-1}, \ldots, t_{m}^{-1}\right] \tag{19.11}
\end{equation*}
$$

such that $\omega\left(\mathfrak{T}_{v}\right)$ has the following expressions:
(I) $v \nmid c$.

$$
\begin{aligned}
& \omega\left(\mathfrak{T}_{v}\right)=\frac{1-[\mathfrak{p}]}{1-q^{n}[\mathfrak{p}]} \prod_{i=1}^{n} \frac{1-q^{2 i}\left[\mathfrak{p}^{2}\right]}{\left(1-q^{n} t_{i}[\mathfrak{p}]\right)\left(1-q^{n} t_{i}^{-1}[\mathfrak{p}]\right)} \quad \text { (Case SP), } \\
& \omega\left(\mathfrak{T}_{v}\right)=\frac{\prod_{i=1}^{2 n}\left(1-(-q)^{i-1}[\mathfrak{q}]\right)}{\prod_{i=1}^{n}\left(1-q^{2 n-2} t_{i}[\mathfrak{q}]\right)\left(1-q^{2 n} t_{i}^{-1}[\mathfrak{q}]\right)} \quad(\text { Case UT, } \mathfrak{p r}=\mathfrak{q}), \\
& \left.\omega\left(\mathfrak{T}_{v}\right)=\frac{\prod_{i=0}^{n-1}\left(1-q^{2 i}\left[\mathfrak{q}^{2}\right]\right)}{\prod_{i=1}^{n}\left(1-q^{n-1} t_{i}[\mathfrak{q}]\right)\left(1-q^{n} t_{i}^{-1}[\mathfrak{q}]\right)} \quad \text { (Case UT, } \mathfrak{p r}=\mathfrak{q}^{2}\right), \\
& \left.\omega\left(\mathfrak{T}_{v}\right)=\prod_{i=1}^{2 n} \frac{1-q^{i-1}\left[\mathfrak{q}_{1} \mathfrak{q}_{2}\right]}{\left(1-q^{2 n} t_{i}^{-1}\left[\mathfrak{q}_{1}\right]\right)\left(1-q^{-1} t_{i}\left[\mathfrak{q}_{2}\right]\right)} \quad \text { (Case UT, } \mathfrak{p r}=\mathfrak{q}_{1} \mathfrak{q}_{2}\right) .
\end{aligned}
$$

(II) $v \mid \boldsymbol{c}$.

$$
\begin{aligned}
& \omega\left(\mathfrak{T}_{v}\right)=\prod_{i=1}^{n}\left(1-q^{n} t_{i}[\mathfrak{p}]\right)^{-1} \quad(\text { Case SP }), \\
& \omega\left(\mathfrak{T}_{v}\right)=\prod_{i=1}^{n}\left(1-q^{k(n-1)} t_{i}[\mathfrak{q}]\right)^{-1} \quad\left(\text { Case UT, } \mathfrak{p}^{k} \mathfrak{r}=\mathfrak{q}^{2}\right) \\
& \omega\left(\mathfrak{T}_{v}\right)=\prod_{i=1}^{n}\left(1-q^{n-1} t_{i}\left[\mathfrak{q}_{1}\right]\right)^{-1}\left(1-q^{n-1} t_{n+i}\left[\mathfrak{q}_{2}\right]\right)^{-1} \quad\left(\text { Case UT, } \mathfrak{p r}=\mathfrak{q}_{1} \mathfrak{q}_{2}\right) .
\end{aligned}
$$

Here $\omega\left(\mathfrak{T}_{v}\right)=\sum_{\xi \in X} \omega\left(C_{v} \xi C_{v}\right)\left[\nu_{\mathfrak{b}}(\xi)\right], X=C_{v} \backslash \mathfrak{X}_{v} / C_{v} ; \mathfrak{p}$ is the prime ideal in $F$ at $v$ and $q=N(\mathfrak{p})$; in Case UT, $\mathfrak{q}$ and $\mathfrak{q}_{i}$ are prime ideals in $K ; \mathfrak{q}_{1} \neq \mathfrak{q}_{2}$.

Proof. For $v \nmid c$ the formulas were given in [S97, Theorem 16.16 and (16.17.5)]; as for the injectivity of $\omega$, see [S97, Proposition 16.14]. Suppose $\mathfrak{p | c}$ and $\mathfrak{p} \nmid e$ in Case UT. We first consider the case $\mathfrak{p r}=\mathfrak{q}^{2 / k}$ with $k=1$ or 2 ; then $N(\mathfrak{q})=q^{k}$. Given a coset $E_{v} d$ with $d \in B_{v}$, we can find an upper triangular $g$ such that $E_{v} d=E_{v} g$; we may assume that the diagonal elements of $g$ are of the forms $\pi^{e_{1}}, \ldots, \pi^{e_{n}}$ with $e_{i} \in \mathbf{Z}$, where $\pi$ is a prime element of $K_{v}$. We then put $\omega_{0}\left(E_{v} d\right)=\prod_{i=1}^{n}\left(q^{-i k} t_{i}\right)^{e_{i}}$. Next, given $C_{v} \sigma C_{v}$ with $\sigma \in Q(e)_{v}$, we take a decomposition $C_{v} \sigma C_{v}=\bigsqcup_{\xi} C_{v} \xi$ with $\xi \in P_{v}$ and put $\omega\left(C_{v} \sigma C_{v}\right)=\sum_{\xi} \omega_{0}\left(E_{v} d_{\xi}\right)$.

We then extend $\omega$ to $\mathfrak{R}\left(C_{v}, \mathfrak{X}_{v}\right) \mathbf{Q}$-linearly. It can easily be verified that $\omega$ is a ring-homomorphism. (This is similar to what was done in [S97, §16.13].) Its injectivity can be proved in the same manner as in [S97, Proposition 16.14]. By Lemma 19.2 (3) we have

$$
\begin{aligned}
\omega\left(\mathfrak{T}_{v}\right) & =\sum_{d \in E_{v} \backslash B_{v}} \omega_{0}\left(E_{v} d\right)\left[S\left(\mathfrak{b}^{-1}\right)_{v}: d^{*} S\left(\mathfrak{b}^{-1}\right)_{v} d\right][\operatorname{det}(d) \mathfrak{r}] \\
& =\sum_{d \in E_{v} \backslash B_{v}} \omega_{0}\left(E_{v} d\right)|\operatorname{det}(d)|^{-n}[\operatorname{det}(d) \mathfrak{r}],
\end{aligned}
$$

since $\left[S\left(\mathfrak{b}^{-1}\right)_{v}: d^{*} S\left(\mathfrak{b}^{-1}\right)_{v} d\right]=|\operatorname{det}(d)|^{-n}$ by [S97, Lemma 13.2], where $|\mid$ is the valuation of $K_{v}$ such that $|\pi|=q^{-k}$. The last sum is essentially the series $\mathcal{B}$ of [S97, Lemma 16.3], and so using the formula for $\mathcal{B}$ given there, we obtain the desired formula for $\omega\left(\mathfrak{T}_{v}\right)$ in the present case. Case SP can be handled in the same manner. (This is actually done in [S94, Theorem 2.9].) Finally consider the case $\mathfrak{p r}=\mathfrak{q}_{1} \mathfrak{q}_{2}$. In this case $G L_{n}\left(K_{v}\right)$ can be identified with $G L_{n}\left(F_{v}\right) \times G L_{n}\left(F_{v}\right)$, and $E_{v}$ with $E_{v}^{0} \times E_{v}^{0}$, where $E_{v}^{0}=G L_{n}\left(\mathfrak{g}_{v}\right)$. Let $E_{v} d=E_{v}^{0} a \times E_{v}^{0} b$ with $d \in G L\left(K_{v}\right)$ and upper triangular matrices $a$ and $b$ in $G L_{n}\left(F_{v}\right)$ whose diagonal elements are $\pi^{e_{1}}, \ldots, \pi^{e_{n}}$ and $\pi^{e_{n+1}}, \ldots, \pi^{e_{2 n}}$, respectively, where $\pi$ is a prime element of $F_{v}$. Putting $\omega_{0}\left(E_{v} d\right)=\prod_{i=1}^{n}\left(q^{-i} t_{i}\right)^{e_{i}}\left(q^{-i} t_{n+i}\right)^{e_{n+i}}$, we define $\omega\left(C_{v} \sigma C_{v}\right)$ in the same manner as in the above case, and repeat the calculation with these modifications to obtain the desired result.
19.9. Lemma. With a fixed $v \in \mathbf{h}$ prime to $\mathfrak{e}$, put $\omega\left(\mathfrak{T}_{v}\right)=\mathcal{T}_{v}\left(t_{1}, \ldots, t_{m}\right)$ with a rational expression $\mathcal{I}_{v}$ defined for each fixed $v$ as in Theorem 19.8. Let $\mathfrak{R}^{v}$ be the subalgebra of $\mathfrak{R}\left(C_{v}, \mathfrak{X}_{v}\right)$ generated over $\mathbf{Q}$ by $T_{v}(\mathfrak{a})$ for all integral $\mathfrak{r}_{v}$-ideals $\mathfrak{a}$, where $T_{v}(\mathfrak{a})$ is the sum of all $C_{v} \xi C_{v}$ with $\xi \in \mathfrak{X}_{v}$ such that $\nu_{\mathfrak{b}}(\xi)=\mathfrak{a}$. Let $\lambda$ be a $\mathbf{Q}$-linear ring-homomorphism of $\mathfrak{R}^{v}$ into $\mathbf{C}$ that maps the identity element to 1. Then there exist $m$ elements $\mu_{i}$ of $\mathbf{C}$ such that $\mu_{i} \neq 0$ if $v \nmid c$ and that the series $\sum_{\mathfrak{a}} \lambda\left(T_{v}(\mathfrak{a})\right)[\mathfrak{a}]$ coincides with the expression $\mathcal{T}_{v}\left(\mu_{1}, \ldots, \mu_{m}\right)$.

Proof. First assume that $K \neq F, v \nmid c$, and $\mathfrak{p r}=\mathfrak{q}$. Let $P(X)$ be the polynomial in an indeterminate $X$ such that $P([\mathfrak{q}])$ coincides with the denominator of $\omega\left(\mathfrak{T}_{v}\right)$ given in Theorem 19.8. Then we see that $P$ has coefficients in $\omega\left(\mathfrak{R}^{v}\right)$ and hence $t_{i}$ and $t_{i}^{-1}$ are integral over $\omega\left(\mathfrak{R}^{v}\right)$. Identify $\mathfrak{R}^{v}$ with $\omega\left(\mathfrak{R}^{v}\right)$. Then the integrality guarantees that $\lambda: \omega\left(\mathfrak{R}^{v}\right) \rightarrow \mathbf{C}$ can be extended to a homomorphism $\lambda^{\prime}$ of $\mathbf{Q}\left[t_{1}, \ldots, t_{n}, t_{1}^{-1}, \ldots, t_{n}^{-1}\right]$ into $\mathbf{C}$. Putting $\mu_{i}=\lambda^{\prime}\left(t_{i}\right)$ and applying $\lambda^{\prime}$ to $\omega\left(\mathfrak{T}_{v}\right)$, we obtain our assertion in the present situation. If $v \nmid c$ and $\mathfrak{p r}=\mathfrak{q}_{1} \mathfrak{q}_{2}$, then the denominator can be written $P_{1}\left(\left[\mathfrak{q}_{1}\right]\right) P_{2}\left(\left[\mathfrak{q}_{2}\right]\right)$ with polynomials $P_{i}$, and we again see that $t_{i}$ and $t_{i}^{-1}$ are integral over $\omega\left(\mathfrak{R}^{v}\right)$, and hence we obtain the same conclusion. All the remaining cases can be handled in the same manner, except that if $v \mid c$, then the quantities $t_{i}^{-1}$ do not appear, so that we cannot say that $\mu_{i} \neq 0$ for $v \mid c$.

We now present a lemma concerning a generalized Möbius function $\mu$ defined on the set of $\mathfrak{r}$-submodules of a torsion $\mathfrak{r}$-module, which will play an essential role in the next section. Here $\mathfrak{r}$ is the ring of algebraic integers in an arbitrary algebraic number field $K$. We denote by $\mathbf{k}$ the set of all nonarchimedean primes of $K$.
19.10. Lemma. To every finitely generated torsion $\mathfrak{r}$-module $A$ we can uniquely assign an integer $\mu(A)$ so that

$$
\sum_{B \subset A} \mu(B)= \begin{cases}1 & \text { if } A=\{0\}  \tag{19.12a}\\ 0 & \text { if } A \neq\{0\}\end{cases}
$$

Moreover $\mu$ has the following properties:

$$
\sum_{B \subset A} \mu(A / B)= \begin{cases}1 & \text { if } A=\{0\}  \tag{19.12b}\\ 0 & \text { if } A \neq\{0\} .\end{cases}
$$

(19.12c) $\mu(A \oplus B)=\mu(A) \mu(B)$ if $\mathfrak{a} A=\mathfrak{b} B=\{0\}$ with relatively prime integral ideals $\mathfrak{a}$ and $\mathfrak{b}$.

$$
\begin{equation*}
\mu\left((\mathfrak{r} / \mathfrak{p})^{r}\right)=(-1)^{r} N(\mathfrak{p})^{r(r-1) / 2} \text { if } 0 \leq r \in \mathbf{Z} \text { and } \mathfrak{p} \text { is a prime ideal in } K . \tag{19.12d}
\end{equation*}
$$

(19.12e) $\mu(A) \neq 0$ if and only if $A$ is annihilated by a squarefree integral ideal.

Proof. We can define $\mu(A)$ inductively by $\mu(A)=-\sum_{B \nsubseteq A} \mu(B)$, starting from $\mu(\{0\})=1$, which shows also the uniqueness. To prove (19.12b), we may assume that $A=L / N$ with two $r$-lattices $L$ and $N$ in $K^{n}$. For every $\mathfrak{r}$-lattice $X$ in $K^{n}$ put $X^{\prime}=\left\{\left.y \in K^{n}\right|^{t} y X \subset \mathfrak{r}\right\}$. Given an $\mathfrak{r}$-submodule $B$ of $A$, take an $\mathfrak{r}$-lattice $M$ so that $N \subset M \subset L$ and $B=M / N$. Put $\psi(B)=\varphi\left(M^{\prime} / L^{\prime}\right)$ with any fixed $\mathfrak{r}$ isomorphism $\varphi$ of $N^{\prime} / L^{\prime}$ onto $A$. Then $\psi$ gives a one-to-one map of the set of all $\mathfrak{r}$-submodules of $A$ onto itself, and $\psi(B) \cong A / B$ and $A / \psi(B) \cong B$. Therefore $\sum_{B \subset A} \mu(A / B)=\sum_{B \subset A} \mu(\psi(B))=\sum_{C \subset A} \mu(C)$, which combined with (19.12a) gives (19.12b). Next, if $A$ and $B$ are as in (19.12c), then every $r$-submodule of $A \oplus B$ is of the form $A^{\prime} \oplus B^{\prime}$ with $\mathfrak{r}$-submodules $A^{\prime}$ of $A$ and $B^{\prime}$ of $B$. Then (19.12c) can be derived from the relation $\mu(A)=-\sum_{C \not \subsetneq_{A}} \mu(C)$ by induction. The formula of (19.12d) follows from the well-known equality $\sum_{r=0}^{n}(-1)^{r} N(\mathfrak{p})^{r(r-1) / 2} c_{r}^{n}=0$ which holds for $n>0$, where $c_{r}^{n}$ denotes the number of $\mathfrak{r}$-submodules of $(\mathfrak{r} / \mathfrak{p})^{n}$ isomorphic to $(\mathfrak{r} / \mathfrak{p})^{r}$. To prove (19.12e), we first observe that $\mu\left(\mathfrak{r} / \mathfrak{p}^{2}\right)=0$ for every prime ideal $\mathfrak{p}$. Given $A$, let $C$ be the maximum $\mathfrak{r}$-submodule of $A$ that is annihilated by a squarefree integral ideal. Suppose $A \neq C$; then $C \neq\{0\}$ and $-\mu(A)=\sum_{D \subset C} \mu(D)+\sum_{B \not \subset C, B \nsubseteq A} \mu(B)$. The first sum on the right-hand side is 0 . Therefore we obtain $\mu(A)=0$ by induction. The converse part follows from (19.12c, d).
19.11. Let $\mathcal{L}$ denote the set of all $\mathfrak{r}$-lattices in $K^{n}$. For $L \in \mathcal{L}$ and $y \in G L_{n}(K)_{\mathbf{A}}$ we denote by $y L$ the $\mathfrak{r}$-lattice in $K^{n}$ such that $(y L)_{v}=y_{v} L_{v}$ for every $v \in \mathbf{k}$. For $L$ and $M$ in $\mathcal{L}$ we define a fractional ideal $\{L / M\}$ and a multiplicative symbol $[L / M]$ (in the sense of $\S 19.5$ ) by

$$
\begin{equation*}
\{L / M\}=\operatorname{det}(y) \mathfrak{r}, \quad[L / M]=[\{L / M\}]=[\operatorname{det}(y) \mathfrak{r}] \tag{19.13}
\end{equation*}
$$

with any $y \in G L_{n}(K)_{\mathbf{A}}$ such that $M=y L$. These are well-defined. Clearly we have $[L / M][M / N]=[L / N]$. If $L, M \in \mathcal{L}$ and $M \subset L$, we can speak of $\mu(L / M)$. Moreover, for each $v \in \mathbf{k}$ we can speak of $\mu\left(L_{v} / M_{v}\right)$ either by viewing $L_{v} / M_{v}$ as an $\mathfrak{r}$-module, or by defining $\mu$ for $\mathfrak{r}_{v}$-modules, which makes no difference. From (19.12c) we easily obtain

$$
\begin{equation*}
\mu(L / M)=\prod_{v \in \mathbf{k}} \mu\left(L_{v} / M_{v}\right) . \tag{19.14}
\end{equation*}
$$

We now take a subset $\Lambda$ of $\mathcal{L}$ satisfying the following condition: if $L \subset H \subset M$, $L, M \in \Lambda$, and $H \in \mathcal{L}$, then $H \in \Lambda$. Fixing an integral ideal $\mathfrak{c}$, we write $L<M$ and $M>L$ if $L \subset M$ and $M_{v}=L_{v}$ for every $v \mid \mathfrak{c}$.
19.12. Lemma. For two functions $\alpha$ and $\beta$ defined on $\Lambda$ with values in a Z-module, we have

$$
\begin{aligned}
& \alpha(L)=\sum_{L<M \in \Lambda} \beta(M) \text { for every } L \in \Lambda \\
& \Longleftrightarrow \beta(L)=\sum_{L<M \in \Lambda} \mu(M / L) \alpha(M) \text { for every } L \in \Lambda, \\
& \alpha(L)=\sum_{L>M \in \Lambda} \beta(M) \text { for every } L \in \Lambda \\
& \Longleftrightarrow \beta(L)=\sum_{L>M \in \Lambda} \mu(L / M) \alpha(M) \text { for every } L \in \Lambda
\end{aligned}
$$

Here and in Lemma 19.14 below each sum may be an infinite sum, and so we have to assume that it is convergent in a suitable sense, or it is a formal sum.

Proof. Assume the first equality. Then, for a fixed $L \in \Lambda$, we have

$$
\begin{aligned}
\sum_{L<M \in \Lambda} \mu(M / L) \alpha(M) & =\sum_{L<M \in \Lambda} \mu(M / L) \sum_{M<H \in \Lambda} \beta(H) \\
& =\sum_{L<H \in \Lambda} \beta(H) \sum_{L<M<H}^{M} \mu(M / L) .
\end{aligned}
$$

The condition $L<M<H$ can be changed into $L \subset M \subset H$. Applying (19.12a) with $H / L$ as $A$ to the last sum, we find that the last double sum equals $\beta(L)$, which proves the first $\Rightarrow$. All the other cases can be proved similarly by employing (19.12a) or (19.12b).

Notice that if $\mathfrak{r}=\mathbf{Z}$, then $n \mapsto \mu(\mathbf{Z} / n \mathbf{Z})$ is the classical Möbius function, and the first half of Lemma 19.12 is exactly the classical Möbius inversion formula if we take $\Lambda=\{n \mathbf{Z} \mid 0<n \in \mathbf{Z}\}$ and consider $\alpha(n \mathbf{Z})$ a function of $n$.
19.13. Lemma. For any fixed $L \in \mathcal{L}$ we have

$$
\begin{gathered}
\sum_{L \subset M \in \mathcal{L}}[M / L]=\sum_{L \supset M \in \mathcal{L}}[L / M]=\sum_{x \in B / E}[\operatorname{det}(x) \mathfrak{r}]=\prod_{i=1}^{n} \prod_{v \in \mathbf{k}}\left(1-q_{v}^{i-1}\left[\mathfrak{q}_{v}\right]\right)^{-1}, \\
\sum_{L \subset M \in \mathcal{L}} \mu(M / L)[M / L]=\sum_{L \supset M \in \mathcal{L}} \mu(L / M)[L / M]=\prod_{i=1}^{n} \prod_{v \in \mathbf{k}}\left(1-q_{v}^{i-1}\left[\mathfrak{q}_{v}\right]\right),
\end{gathered}
$$

where $\mathfrak{q}_{v}$ is the prime ideal at $v$ and $q_{v}=N\left(\mathfrak{q}_{v}\right)$.
Proof. Putting $M={ }^{t} x^{-1} L$ or $M=x L$ with $x \in B$, we see that the first two sums equal to the third sum, which is clealy the product of $\sum_{x \in B_{v} / E_{v}}\left[\operatorname{det}(x) \mathfrak{r}_{v}\right]$ for all $v \in \mathbf{k}$. Each such sum is determined by [S97, Lemma 3.13], and so we obtain the first line of equalities. Notice that the sums are independent of $L$. Next, the product of the first sums of the two lines (for a fixed $L$ ) is

$$
\sum_{L \subset M \in \mathcal{L}} \mu(M / L)[M / L] \sum_{M \subset N \in \mathcal{L}}[N / M]=\sum_{L \subset N \in \mathcal{L}}[N / L] \sum_{L \subset M \subset N} \mu(M / L)
$$

Applying (19.12a) to the last sum, we see that the double sum is 1 . Similarly the product of the second sums of the two lines is 1 . Therefore we obtain the second line of equalities.
19.14. Lemma. Let $\alpha$ and $\gamma$ be functions defined on $\Lambda$ with values in a $\mathbf{Z}$ module $Y$, and let $\delta$ be a function on $\Lambda$ with values in $\operatorname{End}(Y)$. If

$$
\alpha(L)=\sum_{L<N \in \Lambda} \delta(N) \sum_{H \in \Lambda . L+H=N} \gamma(H) \text { for every } L \in \Lambda,
$$

then

$$
\sum_{L<M \in \Lambda} \mu(M / L) \alpha(M)=\sum_{L<M \in \Lambda} \mu(M / L) \delta(M) \sum_{L \supset H \in \Lambda} \gamma(H) \text { for every } L \in \Lambda
$$

Proof. For fixed $L$ and $N$ in $\Lambda$ we have

$$
\sum_{H \in \Lambda, L+H=N} \gamma(H)=\sum_{N \supset H \in \Lambda}\left(\sum_{L+H \subset M \subset N} \mu(N / M)\right) \gamma(H)=\sum_{L \subset M \subset N} \mu(N / M) \varepsilon(M)
$$

where $\varepsilon(M)=\sum_{M \supset H \in \Lambda} \gamma(H)$, and hence

$$
\begin{aligned}
\alpha(L) & =\sum_{L<N \in \Lambda} \delta(N) \sum_{L \subset M \subset N} \mu(N / M) \varepsilon(M) \\
& =\sum_{L<M \in \Lambda} \sum_{M<N \in \Lambda} \mu(N / M) \delta(N) \varepsilon(M)=\sum_{L<M \in \Lambda} \beta(M)
\end{aligned}
$$

where $\beta(M)=\sum_{M<N \in \Lambda} \mu(N / M) \delta(N) \varepsilon(M)$. Therefore Lemma 19.12 gives the desired conclusion.

## 20. Dirichlet series obtained from Hecke eigenvalues and Fourier coefficients

20.1. Let us now take the ideals $\mathfrak{b}, \mathfrak{c}$, and $\mathfrak{e}$ as in $\S 19.1$ and the group $C$ as in (19.1); in addition we take a Hecke character $\psi$ of $K$ such that
(20.1) $\psi_{v}(a)=1$ for every $a \in \mathfrak{r}_{v}^{\times}, v \in \mathbf{h}$, such that $a-1 \in \mathfrak{r}_{v} \mathbf{c}_{v}$.

Let $k$ be an integral or a half-integral weight (see §16.4). We assume that

$$
\begin{equation*}
\mathfrak{b}^{-1} \subset 2 \mathfrak{d}^{-1} \quad \text { and } \quad \mathfrak{b} \mathfrak{c} \subset 2 \mathfrak{d} \quad \text { if } k \text { is half-integral. } \tag{20.2}
\end{equation*}
$$

This means that $C \subset D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]$ and $\mathfrak{c} \subset 4 \mathfrak{g}$ if $k$ is half-integral, and so $j_{\alpha}^{k}$ is meaningful if $\alpha \in \mathrm{pr}^{-1}(C)$. Then we denote by $\mathcal{M}_{k}(C, \psi)$ for integral $k$ (resp. half-integral $k$ ) the set of all functions $\mathbf{f}: G_{\mathbf{A}} \rightarrow \mathbf{C}$ (resp. $\mathbf{f}: M_{\mathbf{A}} \rightarrow \mathbf{C}$ ) satisfying the following two conditions:
(20.3a) $\mathbf{f}(\alpha x w)=\psi_{\mathbf{c}}\left(\operatorname{det}\left(a_{w}\right)\right)^{-1} j_{w}^{k}(\mathbf{i})^{-1} \mathbf{f}(x)$ if $\alpha \in G, w \in C,($ resp. $\operatorname{pr}(w) \in C)$ and $w(\mathbf{i})=\mathbf{i}$, where $\psi_{c}=\prod_{v \mid c} \psi_{v}$.
(20.3b) For every $p \in G_{\mathbf{h}}$ (resp. $p \in M_{\mathbf{A}}$ such that $\operatorname{pr}(p) \in G_{\mathbf{h}}$ ) there exists an element $f_{p}$ of $\mathcal{M}_{k}$, called the $p$-component of $\mathbf{f}$, such that $\mathbf{f}(p y)=$ $\left(f_{p} \|_{k} y\right)(\mathbf{i})$ for every $y \in G_{\mathbf{a}}$ (resp. $y \in M_{\mathbf{A}}$ such that $\operatorname{pr}(y) \in G_{\mathbf{a}}$ ).

Clearly $f_{p}$ is uniquely determined by $\mathbf{f}$ and $p$. In general $a_{w}$ for an arbitrary $w \in C$ may not be an element of $G L_{n}(K)_{\mathbf{A}}$, but $\left(a_{w}\right)_{v} \in G L_{n}\left(\mathfrak{r}_{v}\right)$ for $v \mid \mathfrak{c}$, so that $\psi_{\mathbf{c}}\left(\operatorname{det}\left(a_{w}\right)\right)$ is meaningful (see (1.18)). Also, if $k$ is half-integral, $x \in M_{\mathbf{A}}$, and $\operatorname{pr}(x) \in C$, then $x \in \mathfrak{M}$, so that $j_{x}^{k}(\mathbf{i})$ and $\|_{k} x$ are meaningful (see $\left.\S 16.4\right)$. Put $\Gamma^{p}=G \cap p C p^{-1}$ and

$$
\begin{equation*}
\mathcal{M}_{k}\left(\Gamma^{p}, \psi\right)=\left\{f \in \mathcal{M}_{k} \mid f \|_{k} \gamma=\psi_{\mathrm{c}}\left(\operatorname{det}\left(a\left(p^{-1} \gamma p\right)\right)\right) f \text { for every } \gamma \in \Gamma^{p}\right\} \tag{20.4}
\end{equation*}
$$

Here $a(x)$ is the $a$-block of $x$; we consider only the case $p=1$ if $k \notin \mathbf{Z}^{\mathbf{b}}$; notice that $\Gamma^{1} \subset \Gamma^{\theta}$ by (20.2), so that $f \|_{k} \gamma$ is meaningful for $\gamma \in \Gamma^{1}$. It can easily be verified that $f_{p}$ of (20.3b) belongs to $\mathcal{M}_{k}\left(\Gamma^{p}, \psi\right)$. Now, by [S97, Lemma 8.12] we can take a finite subset $\mathcal{B}$ of $G_{\mathbf{h}}$ so that

$$
\begin{equation*}
G_{\mathbf{A}}=\bigsqcup_{q \in \mathcal{B}} G q C \text { and } q_{v}=1 \text { for every } q \in \mathcal{B} \text { and every } v \mid c . \tag{20.5}
\end{equation*}
$$

Then for integral $k$ we can show that the map $\mathbf{f} \mapsto\left(f_{q}\right)_{q \in \mathcal{B}}$ is a bijection of $\mathcal{M}_{k}(C, \psi)$ onto $\prod_{q \in \mathcal{B}} \mathcal{M}_{k}\left(\Gamma^{q}, \psi\right)$ (see [S97, Lemma 10.8]). In this situation we write $\mathbf{f} \leftrightarrow\left(f_{q}\right)_{q \in \mathcal{B}}$. We call $\mathbf{f}$ a cusp form if $f_{p} \in \mathcal{S}_{k}$ for every $p$, which is the case if $f_{q} \in \mathcal{S}_{k}$ for every $q \in \mathcal{B}$. We denote by $\mathcal{S}_{k}(C, \psi)$ the set of all cusp forms contained in $\mathcal{M}_{k}(C, \psi)$. If $\psi$ is the trivial character, then we denote $\mathcal{M}_{k}(C, \psi)$ and $\mathcal{S}_{k}(C, \psi)$ by $\mathcal{M}_{k}(C)$ and $\mathcal{S}_{k}(C)$. Clearly $\mathcal{M}_{k}(C, \psi)=\mathcal{M}_{k}(C)$ and $\mathcal{S}_{k}(C, \psi)=\mathcal{S}_{k}(C)$ for any $\psi$ if $\mathfrak{e}=\boldsymbol{c}$. Notice that $\mathcal{M}_{k}\left(\Gamma^{p}, \psi\right) \neq\{0\}$ only if
(20.6) $\quad \psi_{\mathbf{a}}\left(\zeta^{n}\right)=\zeta^{k}$ for every root of unity $\zeta \in K$ such that $\zeta-1 \in \mathfrak{r e}$.

In Case SP we have $G_{\mathbf{A}}=G C$ because of strong approximation in $G=S p(n, F)$, and so we can take $\mathcal{B}=\{q\}$ with any $q \in G_{\mathbf{h}}$, and hence $\mathbf{f} \mapsto f_{q}$ is a bijection of $\mathcal{M}_{k}\left(C, \psi_{i}\right)$ onto $\mathcal{M}_{k}\left(\Gamma^{q}, \psi\right)$. The same conclusion holds for half-integral $k$ with $q=1$. Indeed, we first observe that $M_{\mathbf{A}}=G \cdot \operatorname{pr}^{-1}(C)$. Given $f \in \mathcal{M}_{k}\left(\Gamma^{1}, \psi\right)$, define $\mathbf{f}: M_{\mathbf{A}} \rightarrow \mathbf{C}$ by $\mathbf{f}(\alpha w)=\psi_{\mathbf{c}}\left(\operatorname{det}\left(a_{w}\right)\right)^{-1}\left(f \|_{k} w\right)(\mathbf{i})$ for $\alpha \in G$ and $w \in$ $\mathrm{pr}^{-1}(C)$. Then this is well-defined and satisfies (20.3a). Now given $p \in \operatorname{pr}^{-1}\left(G_{\mathbf{h}}\right)$ and $y \in \operatorname{pr}^{-1}\left(G_{\mathbf{a}}\right)$, take $\alpha \in G$ and $x \in \mathrm{pr}^{-1}(C)$ so that $p=\alpha x$. Then $\mathbf{f}(p y)=$ $\psi_{\mathfrak{c}}\left(\operatorname{det}\left(a_{x}\right)\right)^{-1}\left(f \|_{k} x y\right)(\mathbf{i})$. Since $\operatorname{pr}(x)_{\mathbf{a}}=\alpha^{-1}$, if we choose a suitable element $\xi=\left(\alpha^{-1}, t(z)\right)$ in the group $\mathcal{G}$ of $\S 14.14$, then $f\left\|_{k} x=f\right\|_{k} \xi \in \mathcal{M}_{k}$. Thus $\mathbf{f}$ satisfies (20.3b). Clearly $f_{1}=f$. This proves the surjectivity of the map. The injectivity follows immediately from (20.3b).

Now, with $S$ as in (16.1a), put

$$
\begin{equation*}
S_{+}=\left\{\xi \in S \mid \xi_{v} \geq 0 \quad \text { for every } \quad v \in \mathbf{a}\right\} \tag{20.7}
\end{equation*}
$$

20.2. Proposition. Given $\mathbf{f} \in \mathcal{M}_{k}(C, \psi)$, there is a complex number $c(\tau, q ; \mathbf{f})$, written also $c_{\mathbf{f}}(\tau, q)$, determined for $\tau \in S_{+}$and $q \in G L_{n}(K)_{\mathbf{A}}$, such that

$$
\mathbf{f}\left(r_{P}\left[\begin{array}{cc}
q & s \widehat{q}  \tag{20.8}\\
0 & \widehat{q}
\end{array}\right]\right)=\operatorname{det}(q)_{\mathbf{a}}^{[k] \rho}\left|\operatorname{det}(q)_{\mathbf{a}}\right|^{k-[k]} \sum_{\tau \in S_{+}} c(\tau, q ; \mathbf{f}) \mathbf{e}_{\mathbf{a}}^{n}\left(\mathbf{i} q^{*} \tau q\right) \mathbf{e}_{\mathbf{A}}^{n}(\tau s)
$$

for every $s \in S_{\mathbf{A}}$, where $r_{P}$ should be ignored if $k$ is integral, and $[k]$ is the integral part of $k$ as defined in §16.4. Moreover $c_{\mathbf{f}}(\tau, q)$ has the following properties:

$$
\begin{align*}
& c_{\mathbf{f}}(\tau, q) \neq 0 \text { only if } \mathbf{e}_{\mathbf{h}}^{n}\left(q^{*} \tau q s\right)=1 \text { for every } s \in S_{\mathbf{h}}\left(\mathfrak{b}^{-1} \mathfrak{e}\right) ;  \tag{20.9a}\\
& c_{\mathbf{f}}(\tau, q)=c_{\mathbf{f}}\left(\tau, q_{\mathbf{h}}\right) ; \\
& c_{\mathbf{f}}\left(b^{*} \tau b, q\right)=\operatorname{det}(b)^{[k] \rho}|\operatorname{det}(b)|^{k-[k]} c_{\mathbf{f}}(\tau, b q) \text { for every } b \in G L_{n}(K) ;  \tag{20.9c}\\
& \psi_{\mathbf{h}}(\operatorname{det}(e)) c_{\mathbf{f}}(\tau, q e)=c_{\mathbf{f}}(\tau, q) \text { for every } e \in E^{\prime}
\end{align*}
$$

Furthermore, if $\beta \in G \cap \operatorname{diag}[r, \widehat{r}] C p^{-1}$ with $r \in G L_{n}(K)_{\mathbf{h}}$ and $p \in G_{\mathbf{h}}$, then

$$
\begin{equation*}
j^{k}\left(\beta, \beta^{-1} z\right) f_{p}\left(\beta^{-1} z\right)=\psi_{\mathbf{c}}\left(\operatorname{det}\left(a_{\beta p}^{-1} r\right)\right) \sum_{\tau \in S_{+}} c_{\mathbf{f}}(\tau, r) \mathbf{e}_{\mathbf{a}}^{n}(\tau z) \tag{20.9e}
\end{equation*}
$$

where $f_{p}$ is the $p$-component of $\mathbf{f}$ and $a_{\beta p}$ is the $a$-block of $\beta p$. Here we take $p=1$ if $k \notin \mathbf{Z}^{\mathbf{b}}$.

Remark. Taking $\beta=1$ and $p=\operatorname{diag}[r, \widehat{r}]$ in (20.9e), we obtain

$$
\begin{equation*}
f_{p}(z)=\sum_{\tau \in S_{+}} c_{\mathbf{f}}(\tau, r) \mathbf{e}_{\mathbf{a}}^{n}(\tau z) \quad \text { if } \quad p=\operatorname{diag}[r, \widehat{r}] . \tag{20.9f}
\end{equation*}
$$

If $k \notin \mathbf{Z}^{\mathbf{b}}$ and $p=1$, then $\beta$ of (20.9e) belongs to $\mathfrak{M}$, so that $j_{\beta}^{k}$ is meaningful.
Proof. We first consider the case of integral $k$. Let $x=\left[\begin{array}{cc}q & s \widehat{q} \\ 0 & \widehat{q}\end{array}\right]$ with $q \in$ $G L_{n}(K)_{\mathbf{A}}$ and $s \in S_{\mathbf{A}}$; put $p=x_{\mathbf{h}}$, define $f_{p}$ by (20.3b), and put $f_{p}(z)=$
$\sum_{\tau \in S_{+}} c(\tau) \mathbf{e}_{\mathbf{a}}^{n}(\tau z)$. Then $x(\mathbf{i})=i q_{\mathbf{a}} q_{\mathbf{a}}^{*}+s_{\mathbf{a}}$. Since $c(\tau)$ depends only on $p$, we can put

$$
\mathbf{f}(x)=\mathbf{f}\left(p x_{\mathbf{a}}\right)=\left(f_{p} \| x\right)(\mathbf{i})=\operatorname{det}\left(q^{*}\right)^{k} \sum_{\tau \in S_{+}} c(\tau, q, s) \mathbf{e}_{\mathbf{a}}^{n}\left(i q^{*} \tau q\right) \mathbf{e}_{\mathbf{A}}^{n}(\tau s)
$$

with $c(\tau, q, s)=\mathbf{e}_{\mathbf{h}}^{n}(-\tau s) c(\tau)$, which depends only on $q_{\mathbf{h}}, s_{\mathbf{h}}$, and $\tau$. Now $\mathbf{f}(\alpha x w)$ $=\mathbf{f}(x)$ for every $\alpha \in R$ and $w \in R_{\mathrm{h}} \cap C$, so that we easily see that $c(\tau, q, s+h)=$ $c(\tau, q, s)$ for every $h \in S+\prod_{v \in \mathbf{h}} M_{v}$ with an $r$-lattice $M$ in $S$ (depending on $q$ ). Thus $c(\tau, q, s)$ is independent of $s$. Therefore we obtain (20.8) and (20.9b). Now given $\beta$, $r$, and $p$ as in the last part of our proposition, take $y \in P_{\mathrm{a}}$ so that $z=y(\mathbf{i})$. Then we can put $\beta^{-1} \operatorname{diag}[r, \widehat{r}] y=p w$ with $w \in C$. By (20.3a, b), $\mathbf{f}(\operatorname{diag}[r, \widehat{r}] y)=\mathbf{f}(p w)=\psi_{\mathbf{c}}\left(\operatorname{det}\left(a_{w}\right)\right)^{-1} \mathbf{f}\left(p w_{\mathbf{a}}\right)=\psi_{\mathbf{c}}\left(\operatorname{det}\left(a_{w}\right)\right)^{-1}\left(f_{p} \| w\right)(\mathbf{i})$. Since $\left(\beta^{-1} y\right)_{\mathbf{a}}=w_{\mathbf{a}}$, we have $f_{p}\left\|w=f_{p}\right\|\left(\beta^{-1} y\right)$. Observing that $\psi_{c}\left(\operatorname{det}\left(a_{w}\right)\right)=$ $\psi_{c}\left(\operatorname{det}\left(a_{\beta p}^{-1} r\right)\right)$, we have $\left(f_{p} \|\left(\beta^{-1} y\right)\right)(\mathbf{i})=\psi_{c}\left(\operatorname{det}\left(a_{\beta p}^{-1} r\right)\right) \mathbf{f}(\operatorname{diag}[r, \widehat{r}] y)$. Applying (20.8) to $\mathbf{f}(\operatorname{diag}[r, \widehat{r}] y)$, we obtain (20.9e). The remaining properties of $c_{\mathbf{f}}(\tau, q)$ can easily be verified by means of (20.3a).

As for half-integral $k$, identify every element of $P_{\mathbf{A}}$ with its image under $r_{P}$. Then we can repeat the above argument to obtain our assertion in the same manner, except (20.9e); the only necessary modification is that we have to take $|\operatorname{det}(q)|^{k-[k]} \cdot \operatorname{det}(q)^{[k]}$ instead of $\operatorname{det}\left(q^{*}\right)^{k}$, in view of (16.19). To prove (20.9e), we put $\operatorname{diag}[r, \widehat{r}] y=\beta w$ with $y$ as before and $w \in \operatorname{pr}^{-1}(C)$. Then $\beta^{-1} z=w(\mathbf{i})$ and we have again $\mathbf{f}(\operatorname{diag}[r, \widehat{r}] y)=\psi_{\mathbf{c}}\left(\operatorname{det}\left(a_{w}\right)\right)^{-1}\left(f_{1} \| w\right)(\mathbf{i})$. By (16.16c) and (16.19),

$$
\begin{equation*}
j_{\beta}^{k}(w(\mathbf{i})) j_{w}^{k}(\mathbf{i})=j_{\beta w}^{k}(\mathbf{i})=\operatorname{det}\left(d_{y}\right)^{[k]}\left|\operatorname{det}\left(d_{y}\right)\right|^{k-[k]} \tag{*}
\end{equation*}
$$

Now $j_{w}^{k}(\mathbf{i})^{-1} f_{1}\left(\beta^{-1} z\right)=\left(f_{1} \| w\right)(\mathbf{i})$. Therefore, employing (*) and applying (20.8) to $\mathbf{f}(\operatorname{diag}[r, \widehat{r}] y)$, we obtain (20.9e).
20.3. Assuming $k$ to be integral, we now define the action of $\mathfrak{R}(C, \mathfrak{X})$ on $\mathcal{M}_{k}(C, \psi)$; the case of half-integral $k$ will be explained in Section 21. We first make the following observation. If $\xi$ belongs to the set $\mathfrak{X}$ of (19.2c), then $\left(a_{\xi}\right)_{v}$ is invertible for every $v \mid \boldsymbol{c}$, and so $\psi_{\mathbf{c}}\left(\operatorname{det}\left(a_{\xi}\right)\right)$ is meaningful. Put $\varphi(\xi)=\psi_{\mathbf{c}}\left(\operatorname{det}\left(a_{\xi}\right)\right)$. Then $\varphi(\alpha \xi \beta)=\varphi(\alpha) \varphi(\xi) \varphi(\beta)$ for $\alpha, \beta \in C$.

Now, given $\xi \in \mathfrak{X}$ and $\mathbf{f} \in \mathcal{M}_{k}(C, \psi)$, take a finite subset $Y$ of $G_{\mathbf{h}}$ so that $C \xi C=\bigsqcup_{\eta \in Y} C y$ and define $\mathbf{f} \mid C \xi C: G_{\mathbf{A}}^{\varphi} \rightarrow \mathbf{C}$ by

$$
\begin{equation*}
(\mathbf{f} \mid C \xi C)(x)=\sum_{y \in Y} \psi_{c}\left(\operatorname{det}\left(a_{y}\right)\right)^{-1} \mathbf{f}\left(x y^{-1}\right) \quad\left(x \in G_{\mathbf{A}}^{\varphi}\right) \tag{20.10}
\end{equation*}
$$

We can easily verify that this does not depend on the choice of $Y$, and also that $\mathbf{f} \mid C \xi C \in \mathcal{M}_{k}(C, \psi)$. This action can be extended linearly to the whole $\mathfrak{R}(C, \mathfrak{X})$. We easily see that it defines a ring-homomorphism of $\mathfrak{R}(C, \mathfrak{X})$ into $\operatorname{End}\left(\mathcal{M}_{k}(C, \psi)\right)$. We then define a formal Dirichlet series $\mathbf{f} \mid \mathfrak{T}$ with coefficients in $\mathcal{M}_{k}(C, \psi)$ by

$$
\begin{equation*}
\mathbf{f} \mid \mathfrak{T}=\sum_{\xi \in C \backslash \mathfrak{x} / C}(\mathbf{f} \mid C \xi C)\left[\nu_{\mathfrak{b}}(\xi)\right] \tag{20.11}
\end{equation*}
$$

where $\nu_{\mathfrak{b}}(\xi)$ is defined by (19.4). Clearly $\mathbf{f} \mid \mathfrak{T}=\sum_{\mathfrak{a}}(\mathbf{f} \mid T(\mathfrak{a}))[\mathfrak{a}]$. Notice also that $(\mathbf{f} \mid C \xi C)(x)=\sum_{y \in Y} \mathbf{f}\left(x y^{-1}\right)$ if $\mathfrak{e}=\mathfrak{c}$.

As in $\S 19.11$ let $\mathcal{L}$ denote the set of all $\mathfrak{r}$-lattices in $K_{1}^{n}$. We put $L_{0}=\mathfrak{r}_{1}^{n}$ and we shall often express an element $L$ of $\mathcal{L}$ in the form $L=y L_{0}$ with $y \in G L_{n}(K)_{\mathbf{h}}$. Let $\mathfrak{d}, \delta$, and $\widetilde{S}$ be as in $\S 16.1$. For $\tau \in S$ put

$$
\begin{equation*}
\mathcal{L}_{\tau}=\left\{L \in \mathcal{L} \mid \ell^{*} \tau \ell \in \mathfrak{b} \mathfrak{e}^{-1} \mathfrak{d}^{-1} \quad \text { for every } \quad \ell \in L\right\} \tag{20.12}
\end{equation*}
$$

We easily see that $\mathcal{L}_{\tau}$ consists of all the $r$-lattices $y L_{0}$ with $y \in G L_{n}(K)_{\mathbf{h}}$ such that $y^{*} \tau y \in \mathfrak{b e}^{-1} \mathfrak{d}^{-1} \widetilde{S}$. Notice that if $L \in \mathcal{L}$ and $L \subset H \in \mathcal{L}_{\tau}$, then $L \in \mathcal{L}_{\tau}$; moreover, if $\operatorname{det}(\tau) \neq 0$, the set $\left\{M \in \mathcal{L}_{\tau} \mid L \subset M\right\}$ for a fixed $L$ is finite. For $L$ and $M$ in $\mathcal{L}$ let us write $L<M$ if $L \subset M$ and $L_{v}=M_{v}$ for every $v \mid c$.

We consider the Fourier expansion of Proposition 20.2 and investigate their relationship with the formal series $\mathfrak{T}$. By (20.9a), for $y \in G L_{n}(F)_{\mathbf{h}}$ we have

$$
\begin{equation*}
c(\tau, y ; \mathbf{f}) \neq 0 \quad \Longrightarrow \quad y^{*} \tau y \in \mathfrak{b} \mathfrak{e}^{-1} \mathfrak{d}^{-1} \widetilde{S} \quad \Longleftrightarrow \quad y L_{0} \in \mathcal{L}_{\tau} \tag{20.13}
\end{equation*}
$$

Now our first main result of this section can be stated as follows:
20.4. Theorem. Given $\tau \in S_{+}, L \in \mathcal{L}_{\tau}$, and $\mathbf{f} \in \mathcal{M}_{k}(C, \psi)$, take $q \in G L_{n}(K)_{\mathbf{h}}$ so that $L=q L_{0}$ and define formal Dirichlet series $D(\tau, q ; \mathbf{f}), a(\tau, L)$, and $A(\tau, L)$ by

$$
\begin{gathered}
D(\tau, q ; \mathbf{f})=\sum_{x \in B^{\prime} / E^{\prime}} \psi_{\mathfrak{c}}(\operatorname{det}(q x))|\operatorname{det}(x)|_{K}^{-\kappa} c(\tau, q x ; \mathbf{f})[\operatorname{det}(x) \mathfrak{r}], \\
A(\tau, L)=|\operatorname{det}(q)|_{K}^{-\kappa}\left[\operatorname{det}\left(q q^{*}\right) \mathfrak{r}\right] \sum_{L<M \in \mathcal{L}_{\tau}} \mu(M / L) a(\tau, M), \\
a(\tau, L)=|\operatorname{det}(q)|_{K}^{\kappa}\left[\operatorname{det}\left(q q^{*}\right)^{-1} \mathfrak{r}\right] \alpha_{\mathfrak{c}}^{0}\left(\varepsilon_{b} q^{*} \tau q\right) .
\end{gathered}
$$

Here $|w|_{K}$ denotes the idele norm of $w \in K_{\mathbf{A}}^{\times} ; \kappa=n+1$ in Case SP and $\kappa=n$ in Case UT; $\mu(M / L)$ is the Möbius function introduced in Section 19; $\varepsilon_{b}$ is an element of $F_{\mathbf{h}}^{\times}$such that $\varepsilon_{b} \mathfrak{g}=\mathfrak{b}^{-1} \mathfrak{d} ; \alpha_{\mathfrak{c}}^{0}$ is the series of (19.6). Then

$$
\begin{gathered}
{[\operatorname{det}(\widehat{q}) \mathfrak{r}] A(\tau, L) D(\tau, q ; \mathbf{f})=\sum_{L<M \in \mathcal{L}_{\tau}} \mu(M / L) \psi_{\mathfrak{c}}(\operatorname{det}(y))[\operatorname{det}(\widehat{y}) \mathfrak{r}] c(\tau, y ; \mathbf{f} \mid \mathfrak{T}),} \\
{[\operatorname{det}(\widehat{q}) \mathfrak{r}] \psi_{\mathfrak{c}}(\operatorname{det}(q)) c(\tau, q ; \mathbf{f} \mid \mathfrak{T})=\sum_{L<M \in \mathcal{L}_{\tau}}[\operatorname{det}(\widehat{y}) \mathfrak{r}] A(\tau, M) D(\tau, y ; \mathbf{f}),}
\end{gathered}
$$

where $y$ in the last two sums is an element of $G L_{n}(K)_{\mathbf{h}}$ chosen for each $M$ so that $M=y L_{0}$ and $y^{-1} q \in B^{\prime}$. In particular, if $\mathbf{f} \mid T(\mathfrak{a})=\lambda(\mathfrak{a}) \mathbf{f}$ with $\lambda(\mathfrak{a}) \in \mathbf{C}$ for every integral $\mathfrak{r}$-ideal $\mathfrak{a}$, then

$$
\begin{gathered}
A(\tau, L) D(\tau, q ; \mathbf{f})=\sum_{\mathfrak{a}} \lambda(\mathfrak{a})[\mathfrak{a}] \sum_{L<M \in \mathcal{L}_{\tau}} \mu(M / L) \psi_{\mathfrak{c}}(\operatorname{det}(y))\left[\operatorname{det}\left(q^{*} \widehat{y}\right) \mathfrak{r}\right] c_{\mathbf{f}}(\tau, y), \\
\psi_{\mathfrak{c}}(\operatorname{det}(q)) c(\tau, q ; \mathbf{f}) \sum_{\mathfrak{a}} \lambda(\mathfrak{a})[\mathfrak{a}]=\sum_{L<M \in \mathcal{L}_{\tau}}\left[\operatorname{det}\left(q^{*} \widehat{y}\right) \mathfrak{r}\right] A(\tau, M) D(\tau, y ; \mathbf{f}) .
\end{gathered}
$$

Remark. The series $a(\tau, L)$ and $A(\tau, L)$ are defined independently of the choice of $q$, as can be seen from (19.7). $D(\tau, q ; \mathbf{f})$ depends only on ( $\left.\tau, q E^{\prime}, \mathbf{f}\right)$. Thus, if $\mathfrak{e}=\mathfrak{g}$, it is independent of the choice of $q$, and we can put $D(\tau, q ; \mathbf{f})=D(\tau, L ; \mathbf{f})$. Then we can write $D(\tau, M ; \mathbf{f})$ for $D(\tau, y ; \mathbf{f})$ in the above formulas. The sum $\sum_{L<M \in \mathcal{L}_{\tau}}$ is a finite sum if $\operatorname{det}(\tau) \neq 0$. In general it may be an infinite sum. We can show, however, that in all cases $A(\tau, L)$ can be expressed as an easy Euler product times a finite sum which is essentially a lower-dimensional version of $A(\tau, L)$ (see [S94a, Proposition 5.4]).

Proof. Taking a subset $\mathcal{R}$ of $G_{\mathbf{h}}$ that represents $C \backslash \mathfrak{X}$, from (20.10) and (20.11) we obtain

$$
\begin{equation*}
(\mathbf{f} \mid \mathfrak{T})(x)=\sum_{y \in \mathcal{R}} \psi_{\mathbf{c}}\left(\operatorname{det}\left(a_{y}\right)\right)^{-1} \mathbf{f}\left(x y^{-1}\right)\left[\nu_{\mathbf{b}}(y)\right] \quad\left(x \in G_{\mathbf{A}}\right) \tag{20.14}
\end{equation*}
$$

By Lemma 19.4 we can take $\mathcal{R}$ to be the set of all $y=\left[\begin{array}{cc}g^{-1} h & g^{-1} \sigma \widehat{h} \\ 0 & g^{*} \widehat{h}\end{array}\right]$ with
$(g, h) \in E^{\prime} \backslash W /\left(E^{\prime} \times 1\right)$ and $\sigma \in S^{\prime} / g S_{\mathbf{h}}\left(\mathfrak{b}^{-1} \mathfrak{e}\right) g^{*}$. Applying [S97, Proposition 3.9] to $\beta^{-1} y \beta$ with $\beta$ as in (19.5), we obtain $\nu_{\mathfrak{b}}(y)=\operatorname{det}\left(g h^{*}\right) \nu_{0}\left(b_{0} \sigma\right)$. Thus

$$
(\mathbf{f} \mid \mathfrak{T})(x)=\sum \psi_{\mathbf{c}}\left(\operatorname{det}\left(h^{-1} g\right)\right) \mathbf{f}\left(x y^{-1}\right)\left[\operatorname{det}\left(g h^{*}\right) \nu_{0}\left(b_{0} \sigma\right)\right]
$$

Substituting $\left[\begin{array}{cc}q & s \widehat{q} \\ 0 & \widehat{q}\end{array}\right]$ for $x$ and making a straightforward calculation, we find, for every $q \in G L_{n}(K)_{\mathbf{h}}$, that
$c(\tau, q ; \mathbf{f} \mid \mathfrak{T})=\sum_{g, h, \sigma} \psi_{\mathbf{c}}\left(\operatorname{det}\left(h^{-1} g\right)\right) c_{\mathbf{f}}\left(\tau, q h^{-1} g\right) \mathbf{e}_{\mathbf{h}}^{n}\left(-\widehat{h} q^{*} \tau q h^{-1} \sigma\right)\left[\operatorname{det}\left(g h^{*}\right) \nu_{0}\left(b_{0} \sigma\right)\right]$.
By (20.13) we may assume that $g^{*} \widehat{h} q^{*} \tau q h^{-1} g \in \mathfrak{b} e^{-1} \mathfrak{d}^{-1} \widetilde{S}$. Therefore, by Lemma 19.6 we have, with $\varepsilon_{b}$ defined as in our theorem,

$$
\begin{align*}
& c(\tau, q ; \mathbf{f} \mid \mathfrak{T})=\sum_{g, h} \psi_{\mathfrak{c}}\left(\operatorname{det}\left(h^{-1} g\right)\right)|\operatorname{det}(g)|_{K}^{-\kappa}\left[\operatorname{det}\left(g h^{*}\right) \mathfrak{r}\right]  \tag{20.15}\\
& \cdot c\left(\tau, q h^{-1} g ; \mathbf{f}\right) \alpha_{\mathfrak{c}}^{0}\left(\varepsilon_{b} \widehat{h} q^{*} \tau q h^{-1}\right)
\end{align*}
$$

where ( $g, h$ ) runs over $E^{\prime} \backslash W /\left(E^{\prime} \times 1\right)$ under the condition that $\widehat{h} q^{*} \tau q h^{-1} \in$ $\mathfrak{b} \mathfrak{e}^{-1} \mathfrak{d}^{-1} \widetilde{S}$, which is so, by (20.13), if and only if $q h^{-1} L_{0} \in \mathcal{L}_{\tau}$. We now fix $\tau \in S_{+}$ and $p \in G L_{n}(K)_{\mathbf{h}}$, and put

$$
\begin{aligned}
& \Lambda=\left\{M \in \mathcal{L}_{\tau} \mid M_{v}=\left(p L_{0}\right)_{v} \text { for every } v \mid e\right\} \\
& X=\left\{y \in G L_{n}(K)_{\mathbf{h}} \mid y_{v} \in E_{v}^{\prime} \text { for every } v \mid \mathfrak{e}\right\}
\end{aligned}
$$

Now take $L=q L_{0} \in \Lambda$ with $q \in p X$. (If we start from a given $q$, then we define $\Lambda$ with $q$ as $p$.) Then it can easily be seen that $(g, h) \mapsto\left(q h^{-1} L_{0}, q h^{-1} g L_{0}\right)$ gives a one-to-one map of the set of all such ( $g, h$ ) onto the set of all $(N, H)$ in $\Lambda \times \Lambda$ such that $L+H=N$ and $L<N$. Given $M \in \Lambda$, we can choose $y \in p X$ so that $M=y L_{0}$. Then, for a fixed $\tau$ put

$$
\begin{gathered}
c(M)=\psi_{\mathbf{c}}(\operatorname{det}(y)) c(\tau, y ; \mathbf{f})|\operatorname{det}(y)|_{K}^{-\kappa}[\operatorname{det}(y) \mathfrak{r}] \\
c^{\prime}(M)=\psi_{\mathbf{c}}(\operatorname{det}(y)) c(\tau, y ; \mathbf{f} \mid \mathfrak{T})[\operatorname{det}(\widehat{y}) \mathfrak{r}]
\end{gathered}
$$

These are well-defined because of (20.9d). Therefore (20.15) can be written

$$
c^{\prime}(L)=\sum_{L<N \in \Lambda} a(\tau, N) \sum_{L+H=N, H \in \Lambda} c(H)
$$

By Lemma 19.14 we obtain

$$
\sum_{L<M \in \Lambda} \mu(M / L) c^{\prime}(M)=\sum_{L<M \in \Lambda} \mu(M / L) a(\tau, M) \sum_{L \supset H \in \Lambda} c(H)
$$

for every $L \in \Lambda$. Now the condition $L<M \in \Lambda$ is equivalent to $L<M \in \mathcal{L}_{\tau}$, since $M \in \Lambda$ if $L \in \Lambda$ and $L<M$. Therefore we obtain the first equality of our theorem. The second equality follows immediately from this and Lemma 19.12. The last two equalities are immediate consequences of the first two.
20.5. Lemma. Let $\tau \in S_{+} \cap G L_{n}(K)$ and $L=q L_{0} \in \mathcal{L}_{\tau}$; let $\mathbf{b}$ be the set of all primes $v \in \mathbf{h}$ prime to $\mathfrak{c}$ such that $\varepsilon_{b} q^{*} \tau q$ is not regular in the sense of $\S 16.1$. Then

$$
A(\tau, L)= \begin{cases}\prod_{v \in \mathbf{b}} g_{v}([\mathfrak{p}]) \prod_{v \nmid c} h_{v}([\mathfrak{p}])^{-1}(1-[\mathfrak{p}]) \prod_{i=1}^{[n / 2]}\left(1-N(\mathfrak{p})^{2 i}[\mathfrak{p}]^{2}\right) & (\text { Case SP) } \\ \prod_{v \in \mathfrak{b}} g_{v}([\mathfrak{p r}]) \prod_{v \nmid \mathfrak{c}} \prod_{i=1}^{n}\left(1-\left(\theta^{i-1}\right)^{*}(\mathfrak{p}) N(\mathfrak{p})^{i-1}[\mathfrak{p r}]\right) & \text { (Case UT) }\end{cases}
$$

Here $\mathfrak{p}$ is the prime ideal of $F$ at $v$ and $g_{v}$ is a polynomial with constant term 1 and with coefficients in $\mathbf{Z} ; h_{v}=1$ if $n$ is odd, and $h_{v}(t)=1-\rho_{\tau}^{*}(\mathfrak{p}) N(\mathfrak{p})^{n / 2} t$ with the Hecke character $\rho_{\tau}$ of $F$ corresponding to $F\left(c^{1 / 2}\right) / F, c=(-1)^{n / 2} \operatorname{det}(\tau)$, if $n$ is even; $\theta$ is the Hecke character of $F$ corresponding to $K / F$.

Proof. If $L<M \in \mathcal{L}_{\tau}$, then we can put $M=q \widehat{x} L_{0}$ with $x \in B^{\prime}$. Therefore $A(\tau, L)=\prod_{v\rceil c} A_{v}(\tau, L)$ with

$$
\begin{equation*}
A_{v}(\tau, L)=\sum_{x} \mu\left(L_{0} / x^{*} L_{0}\right)|\operatorname{det}(x)|_{K}^{-\kappa}\left[\operatorname{det}\left(x x^{*}\right) \mathfrak{r}\right] \alpha_{v}^{0}\left(x^{-1}\left(\varepsilon_{b} q^{*} \tau q\right)_{v} \widehat{x}\right) \tag{20.16}
\end{equation*}
$$

where $x$ runs over $B_{v} / E_{v}$ under the condition that $x^{-1}\left(\varepsilon_{b} q^{*} \tau q\right)_{v} \widehat{x} \in \widetilde{S}_{v}$. If $v \notin \mathbf{b}$, then $A_{v}(\tau, L)=\alpha_{v}^{0}\left(\left(\varepsilon_{b} q^{*} \tau q\right)_{v}\right)$, which is given by Theorem 16.2 ; if $v \in \mathbf{b}$, then $A_{v}(\tau, L)$ is a finite sum, and each $\alpha_{v}^{0}\left(x^{-1}\left(\varepsilon_{b} q^{*} \tau q\right){ }_{v} \widehat{x}\right)$ is a polynomial times a rational expression given in that theorem. (Notice that, in view of (16.7a), $\alpha_{v}^{0}(\zeta)=$ $A_{\zeta}^{0}([\mathfrak{p r}])$ with $A_{\zeta}^{0}$ of that theorem.) Therefore we obtain our lemma.
20.6. Suppose now $\mathbf{f} \mid T(\mathfrak{a})=\lambda(\mathfrak{a}) \mathbf{f}$ for every $\mathfrak{a}$ as in the last part of Theorem 20.4. We naturally assume that $\mathbf{f} \neq 0$. By Lemma 19.9 , for each $v \in \mathbf{h}$ we can determine complex numbers $\lambda_{v, i}$ so that $\sum_{\mathfrak{a}} \lambda(\mathfrak{a})[\mathfrak{a}]=\prod_{v \in \mathbf{h}} \mathcal{T}_{v}\left(\lambda_{v, 1}, \ldots, \lambda_{v, m}\right)$ with $\mathcal{T}_{v}$ in that lemma. Let $Z_{v}^{-1}$ denote the denominator of the expression for $\mathcal{T}_{v}\left(\lambda_{v, 1}, \ldots, \lambda_{v, m}\right)$ obtained from the expression for $\omega\left(\mathfrak{T}_{v}\right)$ given in Theorem 19.8. Namely, denoting by $\mathfrak{p}$ the prime ideal in $F$ at $v$, we have:
(I) $v \nmid c$.

$$
\begin{aligned}
& Z_{v}=\left(1-N(\mathfrak{p})^{n}[\mathfrak{p}]\right)^{-1} \prod_{i=1}^{n}\left\{\left(1-N(\mathfrak{p})^{n} \lambda_{v, i}[\mathfrak{p}]\right)\left(1-N(\mathfrak{p})^{n} \lambda_{v, i}^{-1}[\mathfrak{p}]\right)\right\}^{-1} \quad(\text { Case SP) }, \\
& \left.Z_{v}=\prod_{i=1}^{n}\left\{\left(1-N(\mathfrak{q})^{n-1} \lambda_{v, i}[\mathfrak{q}]\right)\left(1-N(\mathfrak{q})^{n} \lambda_{v, i}^{-1}[\mathfrak{q}]\right)\right\}^{-1} \quad \text { (Case UT, } \mathfrak{p r}=\mathfrak{q}^{e}\right), \\
& Z_{v}=\prod_{i=1}^{2 n}\left\{\left(1-N\left(\mathfrak{q}_{1}\right)^{2 n} \lambda_{v, i}^{-1}\left[\mathfrak{q}_{1}\right]\right)\left(1-N\left(\mathfrak{q}_{2}\right)^{-1} \lambda_{v, i}\left[\mathfrak{q}_{2}\right]\right)\right\}^{-1} \quad\left(\text { Case UT, } \mathfrak{p r}=\mathfrak{q}_{1} \mathfrak{q}_{2}\right) .
\end{aligned}
$$

$$
\text { (II) } v \mid c, v \nmid e .
$$

$$
Z_{v}=\prod_{i=1}^{n}\left(1-N(\mathfrak{p})^{n} \lambda_{v, i}[\mathfrak{p}]\right)^{-1} \quad(\text { Case SP })
$$

$$
Z_{v}=\prod_{i=1}^{n}\left(1-N(\mathfrak{q})^{n-1} \lambda_{v, i}[\mathfrak{q}]\right)^{-1} \quad\left(\text { Case UT, } \mathfrak{p r}=\mathfrak{q}^{e}\right)
$$

$$
Z_{v}=\prod_{i=1}^{n}\left\{\left(1-N\left(\mathfrak{q}_{1}\right)^{n-1} \lambda_{v, i}\left[\mathfrak{q}_{1}\right]\right)\left(1-N\left(\mathfrak{q}_{2}\right)^{n-1} \lambda_{v, n+i}\left[\mathfrak{q}_{2}\right]\right)\right\}^{-1}
$$

(Case UT, $\left.\mathfrak{p r}=\mathfrak{q}_{1} \mathfrak{q}_{2}\right)$.
(III) $v \mid$ e. $\quad Z_{v}=1 \quad$ (Cases SP and UT).

Then from Theorem 19.8 we obtain

$$
\begin{equation*}
\mathfrak{L} \cdot \sum_{\mathfrak{a}} \lambda(\mathfrak{a})[\mathfrak{a}]=\prod_{v \in \mathbf{h}} Z_{v} \tag{20.17}
\end{equation*}
$$

where $\mathfrak{L}$ is a formal Dirichlet series given by

$$
\begin{aligned}
& \mathfrak{L}=\prod_{\mathfrak{p} \mathfrak{c}}\left\{(1-[\mathfrak{p}]) \prod_{i=1}^{n}\left(1-N(\mathfrak{p})^{2 i}[\mathfrak{p}]^{2}\right)\right\}^{-1} \quad \quad \text { (Case SP) } \\
& \mathfrak{L}=\prod_{\mathfrak{p} \mathfrak{c}} \prod_{i=1}^{2 n}\left(1-\left(\theta^{i-1}\right)^{*}(\mathfrak{p}) N(\mathfrak{p})^{i-1}[\mathfrak{p r}]\right)^{-1} \quad \quad \text { (Case UT), }
\end{aligned}
$$

where $\left(\theta^{i-1}\right)^{*}$ is the same as in Lemma 20.5 , and $\mathfrak{p}$ runs over all the prime $\mathfrak{g}$-ideals prime to $c$. We define also a similar formal Dirichlet series $\mathfrak{L}_{0}$ by

$$
\begin{aligned}
& \mathfrak{L}_{0}=\prod_{\mathfrak{p} \mathfrak{c}} \prod_{i=1}^{[(n+1) / 2]}\left(1-N(\mathfrak{p})^{2 n+2-2 i}[\mathfrak{p}]^{2}\right)^{-1} \quad \text { (Case SP) } \\
& \mathfrak{L}_{0}=\prod_{\mathfrak{p} \nmid c} \prod_{i=1}^{n}\left(1-\left(\theta^{n+i-1}\right)^{*}(\mathfrak{p}) N(\mathfrak{p})^{n+i-1}[\mathfrak{p r}]\right)^{-1} \quad \quad \text { (Case UT) }
\end{aligned}
$$

20.7. Theorem. Let $\mathbf{f}$ and $Z_{v}$ be as above and let $\tau \in S_{+} \cap G L_{n}(K)$ and $L=q L_{0} \in \mathcal{L}_{\tau}$ with $q \in G L_{n}(K)_{\mathbf{h}}$. Then we have

$$
\begin{aligned}
D(\tau, q ; \mathbf{f}) \cdot & \mathfrak{L}_{0} \cdot \prod_{v \in \mathbf{b}} g_{v}([\mathfrak{p r}]) \cdot \prod_{v \nmid c} h_{v}([\mathfrak{p r}])^{-1} \\
& =\prod_{v \in \mathbf{h}} Z_{v} \cdot \sum_{L<M \in \mathcal{L}_{\tau}} \mu(M / L) \psi_{\mathfrak{c}}(\operatorname{det}(y))\left[\operatorname{det}\left(q^{*} \widehat{y}\right) \mathfrak{r}\right] c_{\mathbf{f}}(\tau, y)
\end{aligned}
$$

where $M=y L_{0}$ as in Theorem 20.4, and $\mathbf{b}, g_{v}$, and $h_{v}$ are determined for $\tau$ and $q$ as in Lemma 20.5; we put $h_{v}=1$ in Case UT.

Proof. This follows immediately from (20.17), the third equality of Theorem 20.4 concerning $A(\tau, L) D(\tau, q ; \mathbf{f})$, and Lemma 20.5.
20.8. Lemma. Let $0 \neq \mathrm{f} \in \mathcal{M}_{k}(C, \psi)$ with integral or half-integral $k$ as in §20.1. Then the following three conditions are mutually equivalent:
(1) Case SP: $k_{v} \geq n / 2$ for some $v \in \mathbf{a}$; Case UT: $k_{v}+k_{v \rho} \geq n$ for some $v \in \mathbf{a}$.
(2) Case SP: $k_{v} \geq n / 2$ for every $v \in \mathbf{a}$; Case UT: $k_{v}+k_{v \rho} \geq n$ for every $v \in \mathbf{a}$.
(3) $c_{\mathbf{f}}(\tau, r) \neq 0$ for some $\tau \in S_{+} \cap G L_{n}(K)$ and some $r \in G L_{n}(K)_{\mathbf{h}}$.

Moreover, these conditions are satisfied if $\mathbf{f}$ is a cusp form.
Proof. We first note that $G_{\mathbf{A}}=\bigsqcup_{r \in\{r\}} G \cdot \operatorname{diag}[r, \widehat{r}] C$ with a finite subset $\{r\}$ of $G L_{n}(K)_{\mathbf{h}}$. This is trivial in Case SP since $G_{\mathbf{A}}=G C$; in Case UT the fact is included in [S97, Lemma 9.8 (3)]. Given $\mathbf{f} \neq 0$, take $p \in G_{\mathbf{h}}$ so that the $p$ component of $\mathbf{f}$ is nonzero; we may assume that $p=\operatorname{diag}[r, \widehat{r}]$ with $r \in\{r\}$. Then, in view of (20.9f), the mutual equivalence of (1), (2), and (3) follow from Proposition 6.16. The last assertion follows from (6.42).
20.9. Theorem. Let $0 \neq \mathbf{f} \in \mathcal{M}_{k}(C, \psi)$ with integral $k$; suppose that the conditions of Lemma 20.8 are satisfied. Then there exist $\tau \in S_{+} \cap G L_{n}(K)$ and $r \in G L_{n}(K)_{\mathbf{h}}$ such that

$$
\begin{equation*}
0 \neq \psi_{c}(\operatorname{det}(r)) c(\tau, r ; \mathbf{f} \mid \mathfrak{T})=A\left(\tau, r L_{0}\right) D(\tau, r ; \mathbf{f}) \tag{20.18}
\end{equation*}
$$

Suppose in particular that $\mathbf{f} \mid T(\mathfrak{a})=\lambda(\mathfrak{a}) \mathbf{f}$ as above. Then $c_{\mathbf{f}}(\tau, r) \neq 0$ and

$$
\begin{equation*}
\psi_{\mathfrak{c}}(\operatorname{det}(r)) c_{\mathbf{f}}(\tau, r) \prod_{v \in \mathbf{h}} Z_{v}=D(\tau, r ; \mathbf{f}) \cdot \mathfrak{L}_{0} \cdot \prod_{v \nmid \mathfrak{c}} h_{v}([\mathfrak{p r}])^{-1} \cdot \prod_{v \in \mathbf{b}} g_{v}([\mathfrak{p r}]) \tag{20.19}
\end{equation*}
$$

with the symbols as in Lemma 20.5 and Theorem 20.7.
Proof. We have $c_{\mathbf{f}}(\tau, q) \neq 0$ for some $\tau \in S_{+} \cap G L_{n}(K)$ and some $q \in$ $G L_{n}(K)_{\mathbf{h}}$. Then $D(\tau, q ; \mathbf{f}) \neq 0$. Let $\mathcal{L}_{\tau}^{\prime}$ be the set of all lattices $M \in \mathcal{L}_{\tau}$ such that $q L_{0}<M$ and $M=r L_{0}$ with some $r \in G L_{n}(K)_{\mathbf{h}}$ for which $D(\tau, r ; \mathbf{f}) \neq 0$. Since $\mathcal{L}_{\tau}^{\prime}$ is a finite set containing $q L_{0}$, it has a maximal element. Writing it $r L_{0}$ with $r$ such that $D(\tau, r ; \mathbf{f}) \neq 0$, we obtain (20.18) and (20.19) from Theorems 20.4 and 20.7 .
20.10. Lemma. Let $\mathcal{B}$ be as in (20.5); let $\mathbf{f} \in \mathcal{M}_{k}(C, \psi)$; also let $\mathbf{f} \leftrightarrow\left(f_{p}\right)_{p \in \mathcal{B}}$ and $\mathbf{f} \mid \mathfrak{T} \leftrightarrow\left(g_{q}\right)_{q \in \mathcal{B}}$ in the sense of $\S 20.1$. Further let $R^{p}$ be a complete set of representatives for $\Gamma^{p} \backslash(\mathfrak{X} \cap G)$. Then

$$
g_{q}=\sum_{p \in \mathcal{B}} \sum_{\gamma \in R^{p}} \psi_{\mathbf{c}}\left(\operatorname{det}\left(a\left(p^{-1} \gamma q\right)\right)\right)^{-1}\left(f_{p} \| \gamma\right)\left[\nu_{\mathfrak{b}}\left(p^{-1} \gamma q\right)\right] .
$$

This was given in [S97, (11.9.1) and (11.11.3)] for forms on a unitary group of a general type. The proof given in [S97, Lemma $11.8, \S \S 11.9$ and 11.11] is applicable to the present case. In Case SP the matter is simpler, since we can take $\mathcal{B}=\{1\}$.
20.11. We now assume our $f$ to be a cusp form. Thus, given an eigenform $\mathbf{f} \in \mathcal{S}_{k}(C, \psi)$ as in $\S 20.6$ and a Hecke character $\chi$ of $K$, we put

$$
\begin{align*}
\Lambda_{\mathfrak{c}}^{m}(s, \chi) & = \begin{cases}L_{\mathfrak{c}}(2 s, \chi) \prod_{i=1}^{[m / 2]} L_{\mathfrak{c h}}\left(4 s-2 i, \chi^{2}\right) & (\text { Case SP) }, \\
\prod_{i=1}^{m} L_{\mathfrak{c h}}\left(2 s-i+1, \chi_{1} \theta^{i-1}\right) & (\text { Case UT) }, \\
\mathcal{Z}(s, \mathbf{f}, \chi)=\prod_{v \in \mathbf{h}, v \nmid \mathfrak{h}} Z_{v}\left(\chi^{*}(\mathfrak{q}) N(\mathfrak{q})^{-s}\right), \\
\mathfrak{T}(s, \mathbf{f}, \chi)=\sum_{\mathfrak{a}+\mathfrak{r h}=\mathfrak{r}} \chi^{*}(\mathfrak{a}) \lambda(\mathfrak{a}) N(\mathfrak{a})^{-s}\end{cases} \tag{20.20}
\end{align*}
$$

Here $L_{\mathfrak{c}}$ is defined by (16.9), $\chi_{1}$ is the restriction of $\chi$ to $F_{\mathbf{A}}^{\times}, \mathfrak{h}=\mathfrak{g} \cap$ (the conductor of $\chi), \theta$ is the Hecke character of $F$ corresponding to $K / F$, and $Z_{v}\left(\chi^{*}(\mathfrak{q}) N(\mathfrak{q})^{-s}\right)$ should be understood as follows: In Case SP, it is the expression obtained from $Z_{v}$ of $\S 20.6$ by substituting $\chi^{*}(\mathfrak{p}) N(\mathfrak{p})^{-s}$ for $[\mathfrak{p}]$; similarly in Case UT, it is obtained from $Z_{v}$ by substituting $\left(\chi^{*}\left(\mathfrak{q}_{1}\right) N\left(\mathfrak{q}_{1}\right)^{-s}, \chi^{*}\left(\mathfrak{q}_{2}\right) N\left(\mathfrak{q}_{2}\right)^{-s}\right)$ or $\chi^{*}(\mathfrak{q}) N(\mathfrak{q})^{-s}$ for $\left(\left[\mathfrak{q}_{1}\right],\left[\mathfrak{q}_{2}\right]\right)$ or [ $\mathfrak{q}]$ according as $\mathfrak{p r}=\mathfrak{q}_{1} \mathfrak{q}_{2}$ or $\mathfrak{p r}=\mathfrak{q}^{e}$. From (20.17) we obtain, at least formally,

$$
\begin{equation*}
\mathcal{Z}(s, \mathbf{f}, \chi)=\Lambda_{c}^{2 n}(s / u, \chi) \mathfrak{T}(s, \mathbf{f}, \chi), \quad u=2 /[K: F] . \tag{20.23}
\end{equation*}
$$

20.12. Lemma. (1) The series of (20.22) and the product of (20.21) are absolutely convergent for $\operatorname{Re}(s)>2 n+1$ in Case $S P$ and $\operatorname{Re}(s)>2 n$ in Case UT. (This is preliminary to the stronger result given in the following theorem.)
(2) If $\psi$ is of finite order and the conditions of Lemma 20.8 are satisfied, then the eigenvalues $\lambda(\mathfrak{a})$ generate a finite algebraic extension of $\mathbf{Q}$ that depends on $\mathbf{f}$. (As to the nature of this extension, see Lemma 23.15 below.)
(3) If $\mathfrak{e}=\mathfrak{c}$, the space $\mathcal{S}_{h}(C)$ is spanned by eigenfunctions $\mathbf{f}$ of the above type.

Proof. As explained in [S97, §20.13], the convergence can be reduced to that of $\sum_{\tau \in C \backslash \mathfrak{x}}\left|N\left(\mu_{\mathfrak{b}}(\tau)\right)^{-s}\right|$. Also, as explained in the proof of [S97, Lemma 20.11], $\sum_{\tau \in C \backslash \mathfrak{X}} N\left(\mu_{\mathfrak{b}}(\tau)\right)^{-s}$ has an Euler product, whose Euler factor in Case SP can be
obtained by substituting $\left(N(\mathfrak{p})^{-s}, N(\mathfrak{p})^{i}\right)$ for $\left([\mathfrak{p}], t_{i}\right)$ in the expression for $\omega\left(\mathfrak{T}_{v}\right)$ in Theorem 19.8. Then we find that the series is convergent for $\operatorname{Re}(s)>2 n+1$ in Case SP. Similarly, in Case UT, it is convergent for $\operatorname{Re}(s)>2 n$, as already noted in [S97, Proposition 20.4 (3)]. To prove (2), assuming that $\psi$ is of finite order and the conditions of Lemma 20.8 are satisfied, denote by $D$ the field generated over $\mathbf{Q}$ by the values of $\psi$ and the conjugates of $K$; let $\sigma \in \operatorname{Aut}(\mathbf{C} / D)$. In Lemma 23.14 below we shall establish a nonzero element $\mathbf{f}^{\sigma}$ of $\mathcal{M}_{k}(C, \psi)$ such that $\mathbf{f}^{\sigma} \mid T(\mathfrak{a})=\lambda(\mathfrak{a})^{\sigma} \mathbf{f}^{\sigma}$ for every $\mathfrak{a}$. Since $\mathcal{M}_{k}(C, \psi)$ is finite-dimensional, the $\lambda(\mathfrak{a})$ must belong to a finite algebraic extension of $D$. As for (3), with $\mathfrak{X}$ as in (19.2c) we see that $\tau^{-1} \in \mathfrak{X}$ for every $\tau \in \mathfrak{X}$ if $\mathfrak{e}=\boldsymbol{c}$. Since $\mathfrak{R}(C, \mathfrak{X})$ is commutative, [S97, Proposition 11.7] shows that $C \tau C$ for every $\tau \in \mathfrak{X}$ defines a normal operator on $\mathcal{S}_{h}(C)$. Thus we obtain (3).

We can now state our main theorems about the above Euler product.
20.13. Theorem. The function $\mathcal{Z}(s, \mathbf{f}, \chi)$ can be continued to a meromorphic function on the whole s-plane. Moreover, the Euler product on the right-hand side of (20.21) is convergent, and consequently $\mathcal{Z}(s, \mathbf{f}, \chi) \neq 0$, at least for $\operatorname{Re}(s)>$ $(3 n / 2)+2-[K: F]$.

The proof will be completed in $\S 22.9$. We shall also show in Theorem 22.11 that the bound $(3 n / 2)+2-[K: F]$ is best possible in general.

If $\mathfrak{e}=\mathfrak{c}$, we can state analytic properties of $\mathcal{Z}$ in a better form; for some technical reasons we denote the weight of $\mathbf{f}$ by $h$ instead of $k$.
20.14. Theorem. Let $\mathbf{f}$ be an eigenform contained in $\mathcal{S}_{h}(C)$ with $C$ as above; suppose that $\mathbf{e}=\mathbf{c}$ and $\chi_{\mathbf{a}}(x)=x_{\mathbf{a}}^{\ell}\left|x_{\mathbf{a}}\right|^{i \kappa-\ell}$ with $\ell \in \mathbf{Z}^{\mathbf{a}}$ and $\kappa \in \mathbf{R}^{\mathbf{a}}, \sum_{v \in \mathbf{a}} \kappa_{v}=0$; suppose also that $0 \leq h_{v}-\ell_{v} \leq 1$ for every $v \in \mathbf{a}$ in Case SP. Put

$$
\mathcal{R}(s, \mathbf{f}, \chi)=\mathcal{Z}(s, \mathbf{f}, \chi) \prod_{v \in \mathbf{a}} \Gamma_{v}^{h, \ell}\left((s / u)+\left(i \kappa_{v} / 2\right)\right)
$$

with $u=2 /[K: F]$ and $\Gamma_{v}^{h, \ell}$ defined as follows:
Case SP. $\quad \Gamma_{v}^{h, \ell}(s)=\Gamma_{n}^{1}\left(s+h_{v}-\left(\ell_{v}+n+1\right) / 2\right) g^{n}\left(s, \ell_{v}\right)$, $g^{n}(s, a)= \begin{cases}\Gamma_{n}^{1}\left(s+\frac{a-n}{2}\right) \Gamma\left(s+\frac{a}{2}-\left[\frac{a+n}{2}\right]\right) \quad \text { if } a \geq n, \\ \Gamma_{2 a+2-n}^{1}\left(s+\frac{a-n}{2}\right) \Gamma\left(s-\frac{a}{2}\right) \prod_{i=a+2}^{n} \Gamma(2 s-i) \quad \text { if }(n-2) / 2 \leq a<n ;\end{cases}$
Case UT. $\quad \Gamma_{v}^{h, \ell}(s)=p_{v}(s) q^{2 n}\left(s,\left|2 h_{v \rho}+\ell_{v}\right|\right)$

$$
\begin{gathered}
\cdot \Gamma_{n}^{2}\left(s-n+\frac{h_{v}+h_{v \rho}+\left|d_{v}\right|}{2}\right) \Gamma_{n}^{2}\left(s-n+\frac{\left|2 h_{v \rho}+\ell_{v}\right|}{2}\right), \\
p_{v}(s)= \begin{cases}\Gamma_{n}^{2}\left(s+\frac{\left|2 h_{v \rho}+\ell_{v}\right|}{2}\right) \Gamma_{n}^{2}\left(s+\frac{2 h_{v \rho}+\ell_{v}}{2}\right)^{-1} \quad \text { if } d_{v} \geq 0 \\
\Gamma_{n}^{2}\left(s-\frac{2 h_{v}-\ell_{v}}{2}\right) \Gamma_{n}^{2}\left(s-\frac{2 h_{v \rho}+\ell_{v}}{2}\right)^{-1} & \text { if } d_{v}<0\end{cases} \\
q^{t}(s, a)=\prod_{i=1}^{t-a-1} \Gamma\left(s-\frac{a}{2}-\left[\frac{i}{2}\right]\right) \Gamma\left(s-\frac{a}{2}-i\right)^{-1} \quad(0<t \in \mathbf{Z})
\end{gathered}
$$

Here $\Gamma_{n}^{\iota}$ is defined by (16.47); $d_{v}=h_{v}-h_{v \rho}-\ell_{v}$ for $v \in \mathbf{a}$ in Case UT. Then $\mathcal{R}(s, \mathbf{f}, \chi)$ can be continued to the whole $s$-plane as a meromorphic function with
finitely many poles, which are all simple. The set of poles of $\mathcal{R}(u s, \mathbf{f}, \chi)$ is contained in the set of poles of the function $\mathcal{P}(s)$ defined as follows: in Case SP, $\mathcal{P}$ is the product of Theorem 16.11 defined with $\{2 n, \ell, \chi\}$ as $\{n, k, \chi\}$ there; in Case $U T, \mathcal{P}$ is the product given in [S97, Theorem 19.3] defined with $\left\{2 n, 2 h_{v \rho}+\ell_{v}, \chi\right\}$ as $\left\{n, k_{v}, \chi\right\}$ there.

Some more precise results concerning the poles of $\mathcal{Z}$ and $\mathcal{R}$ in Case SP are given in [S96, Theorems B1 and B2]. The proof of the above theorem will be completed in Section 25. Notice that $p_{v}(s)$ and $q^{t}(s, \ell)$ are polynomials in $s$; in particular, $p_{v}=1$ if $0 \leq d_{v} \leq h_{v}+h_{v \rho}$ and $q^{t}(s, a)=1$ if $a \geq t-1$. By (6.42) we may assume that $h_{v} \geq n / 2$ for every $v \in \mathbf{a}$ in Case SP, which is why $g^{n}(s, a)$ is defined only for $a \geq(n-2) / 2$. (Correction to [S97, Theorem 19.3], lines $3 \sim 4$ from the bottom: Read "the set described in Case I" for "the set of (19.3.1).")

## 21. The Euler products for the forms of half-integral weight

21.1. Let us now briefly indicate that the analogues of the theorems of Section 20 can be proved for the forms of half-integral weight. We content ourselves only with giving definitions and making statements without proofs, since such require lengthy calculations; the reade is referred to [S95b] for details. We first put

$$
\begin{gather*}
U=\left\{\alpha \in M_{\mathbf{A}} \mid \operatorname{pr}(\alpha) \in D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]\right\}  \tag{21.1}\\
\mathcal{Z}=\left\{\alpha \in M_{\mathbf{A}} \mid \operatorname{pr}(\alpha) \in \mathfrak{X}_{0}\right\}, \quad \mathfrak{X}_{0}=D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right] Q(\mathfrak{g}) D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right] \tag{21.2}
\end{gather*}
$$

where $Q(\mathfrak{g})$ is defined by (19.2c). Let $k$ be a half-integral weight. Given an element $\alpha=\xi_{1} \sigma \xi_{2} \in \mathcal{Z}$ with $\xi_{i} \in U$ and $\operatorname{pr}(\sigma) \in \mathfrak{X}_{0}$, we put

$$
\begin{equation*}
J^{k}(\alpha, z)=j^{k}\left(\xi_{1} \xi_{2}, z\right) \quad(z \in \mathcal{H}) \tag{21.3}
\end{equation*}
$$

with $j^{k}$ of (16.17).
21.2. Lemma. (1) $J^{k}(\alpha, z)$ is well-defined.
(2) $J^{k}(\xi, z)=j^{k}(\xi, z) \quad$ if $\xi \in U$.
(3) $J^{k}(\xi \alpha \eta, z)=J^{k}(\xi, \alpha \eta z) J^{k}(\alpha, \eta z) J^{k}(\eta, z)$ if $\alpha \in \mathcal{Z}$ and $\xi, \eta \in U$.
(4) $J^{k}(\alpha, z)=j^{[k]}(\operatorname{pr}(\alpha), z) J^{k-[k]}(\alpha, z)$.

The proof of (1) requires a nontrivial fact [S95b, Lemma 2.2]; see the first paragraph of [S95b, p.32]. Once this is established, the remaining assertions follow immediately from (21.3).
21.3. We now take $C, B^{\prime}$, and $E^{\prime}$ as in $\S 19.1$, and put $\Gamma=G \cap C$. This is $\Gamma^{1}$ of $\S 20.1$, and $\mathcal{M}_{k}(\Gamma, \psi)$ is meaningful. As noted there, $\mathcal{M}_{k}(C, \psi)$ is isomorphic to $\mathcal{M}_{k}(\Gamma, \psi)$. We put $\mathbf{f}=f_{\mathbf{A}}$ if an element $\mathbf{f}$ of $\mathcal{M}_{k}(C, \psi)$ corresponds to $f \in$ $\mathcal{M}_{k}(\Gamma, \psi)$. Given $q \in B^{\prime}$, we take a decomposition $G \cap(C \operatorname{diag}[\widehat{q}, q] C)=\bigsqcup_{\alpha \in R} \Gamma \alpha$ with $R \subset G \cap \mathfrak{X}_{0}$. Then we define $f \mid T_{q}$ and $\mathbf{f} \mid T_{q}$ by

$$
\begin{equation*}
\mathbf{f} \mid T_{q}=\left(f \mid T_{q}\right)_{\mathbf{A}}, \quad\left(f \mid T_{q}\right)(z)=\sum_{\alpha \in R} \psi_{\mathbf{c}}\left(\operatorname{det}\left(a_{\alpha}\right)\right)^{-1} J^{k}(\alpha, z)^{-1} f(\alpha z) \tag{21.4}
\end{equation*}
$$

It can easily be seen that $f \mid T_{q}$ is well-defined and belongs to $\mathcal{M}_{k}(\Gamma, \psi)$. For an integral $\mathfrak{g}$-ideal $\mathfrak{a}$ we denote by $T(\mathfrak{a})$ the sum of $T_{q}$ for all $E^{\prime} q E^{\prime}$ such that $\operatorname{det}(q) \mathfrak{g}=$ $\mathfrak{a}$, and put $\mathbf{f} \mid \mathfrak{T}=\sum_{\mathfrak{a}}(\mathbf{f} \mid T(\mathfrak{a}))[\mathfrak{a}]$. By [S95b, Theorem 4.4] we can establish the analogues of $\mathfrak{R}\left(C_{v}, \mathfrak{X}_{v}\right)$ and $\omega_{v}$ in the present case so that $\omega\left(\mathfrak{T}_{v}\right)=1$ for $v \mid \mathfrak{e}$ and

$$
\begin{array}{ll}
\omega\left(\mathfrak{T}_{v}\right)=\prod_{i=1}^{n} \frac{1-q^{2 i-1}\left[\mathfrak{p}^{2}\right]}{\left(1-q^{n} t_{i}[\mathfrak{p}]\right)\left(1-q^{n} t_{i}^{-1}[\mathfrak{p}]\right)} & (v \nmid \mathfrak{c}) \\
\omega\left(\mathfrak{T}_{v}\right)=\prod_{i=1}^{n}\left(1-q^{n} t_{i}[\mathfrak{p}]\right)^{-1} & \left(v \mid \mathfrak{e}^{-1} \mathfrak{c}\right) \tag{21.5b}
\end{array}
$$

where $q=N(\mathfrak{p})$ and $\mathfrak{p}$ is the prime ideal at $v$. Assuming $\mathbf{f} \mid T(\mathfrak{a})=\lambda(\mathfrak{a}) \mathbf{f}$ for every $\mathfrak{a}$, we have complex numbers $\lambda_{v, i}$ with which we have

$$
\begin{gather*}
\mathfrak{L} \cdot \sum_{\mathfrak{a}} \lambda(\mathfrak{a})[\mathfrak{a}]=\prod_{v \in \mathfrak{h}} Z_{v}  \tag{21.6a}\\
\mathfrak{L}=\prod_{\mathfrak{p} \mathfrak{c}} \prod_{i=1}^{n}\left(1-N(\mathfrak{p})^{2 i-1}[\mathfrak{p}]^{2}\right)^{-1}  \tag{21.6b}\\
Z_{v}= \begin{cases}1 & (v \mid \mathfrak{e}), \\
\prod_{i=1}^{n}\left(1-N(\mathfrak{p})^{n} \lambda_{v, i}[\mathfrak{p}]\right)^{-1}, & \left(v \mid \mathfrak{e}^{-1} \mathfrak{c}\right) \\
\prod_{i=1}^{n}\left\{\left(1-N(\mathfrak{p})^{n} \lambda_{v, i}[\mathfrak{p}]\right)\left(1-N(\mathfrak{p})^{n} \lambda_{v, i}^{-1}[\mathfrak{p}]\right)\right\}^{-1} & (v \nmid \mathfrak{c})\end{cases} \tag{21.6c}
\end{gather*}
$$

where $\mathfrak{p}$ is the prime ideal at $v$.
Given a Hecke character $\chi$ of $F$ of conductor $\mathfrak{h}$, we put

$$
\begin{equation*}
\mathcal{Z}(s, \mathbf{f}, \chi)=\prod_{v \in \mathbf{h}, v \nmid \mathfrak{h}} Z_{v}\left(\chi^{*}(\mathfrak{p}) N(\mathfrak{p})^{-s}\right), \tag{21.7}
\end{equation*}
$$

where $Z_{v}\left(\chi^{*}(\mathfrak{p}) N(\mathfrak{p})^{-s}\right)$ is the expression obtained from $Z_{v}$ by substituting $\chi^{*}(\mathfrak{p})$ $N(\mathfrak{p})^{-s}$ for $[\mathfrak{p}]$.
21.4. Theorem. The formulas and assertions of Theorems 20.4, 20.7, 20.9, and 20.13 and Lemmas 20.5 and 20.12 are true for $\mathbf{f} \in \mathcal{M}_{k}(C, \psi)$ with half-integal $k$ if we make the following modifications: $\mathfrak{L}_{0}=\prod_{\mathfrak{p} \nmid c} \prod_{i=1}^{[n / 2]}\left[1-N(\mathfrak{p})^{2 n+1-2 i}[\mathfrak{p}]^{2}\right]^{-1}$; $\alpha_{\mathrm{c}}^{0}\left(\varepsilon_{b} q^{*} \tau q\right)$ should be replaced by $\left.\alpha_{\mathrm{c}}^{1}{ }^{( }{ }^{t} q \tau q\right)$, with $\alpha_{c}^{1}$ defined by

$$
\alpha_{\mathfrak{c}}^{1}(\zeta)=\prod_{v \nmid \mathfrak{c}} \alpha_{v}^{1}(\zeta, s), \quad \alpha_{v}^{1}(\zeta)=\sum_{\sigma \in S_{v} / S(\mathfrak{g})_{v}} \omega\left(\delta_{v}^{-1} \sigma\right) \mathbf{e}_{v}^{n}\left(-\delta_{v}^{-1} \zeta_{v} \sigma\right)\left[\nu_{0}(\sigma)\right] \quad(\zeta \in \widetilde{S})
$$

where $\omega\left(\delta_{v}^{-1} \sigma\right)$ is as in (16.7b);

$$
A(\tau, L)=\prod_{v \in \mathbf{b}} g_{v}\left(N(\mathfrak{p})^{-1 / 2}[\mathfrak{p}]\right) \prod_{v \nmid \mathfrak{c}}\left\{h_{v}([\mathfrak{p}])^{-1} \prod_{i=1}^{[(n+1) / 2]}\left(1-N(\mathfrak{p})^{2 i-1}[\mathfrak{p}]^{2}\right)\right\}
$$

$\mathbf{b}$ is the set of all $v \in \mathbf{h}$ such that $v \nmid c$ and $\left({ }^{t} q \tau q\right)_{v} \notin E_{v} ; h_{v}=1$ if $n$ is even, and $h_{v}(t)=1-\rho_{\tau}^{*}(\mathfrak{p}) N(\mathfrak{p})^{n / 2} t$ with the Hecke ideal character $\rho_{\tau}^{*}$ of $F$ corresponding to $F\left(c^{1 / 2}\right) / F, c=(-1)^{(n-1) / 2} \operatorname{det}(2 \tau)$, if $n$ is odd; in Lemma 20.12 (2), $\lambda(\mathfrak{a})$ must be replaced by $N(\mathfrak{a})^{1 / 2} \lambda(\mathfrak{a})$.

We add a remark about the last point. Put $T^{\prime}(\mathfrak{a})=N(\mathfrak{a})^{1 / 2} T(\mathfrak{a})$ and $\lambda^{\prime}(\mathfrak{a})=$ $N(\mathfrak{a})^{1 / 2} \lambda(\mathfrak{a})$; take $\sigma \in \operatorname{Aut}(\mathbf{C} / D)$ with $D$ of the proof of Lemma 20.12. Then in Lemma 23.15 below we find an element $\mathbf{f}^{\sigma}$ of $\mathcal{M}_{k}(C, \psi)$ such that $\mathbf{f}^{\sigma} \mid T^{\prime}(\mathfrak{a})=$ $\lambda^{\prime}(\mathfrak{a})^{\sigma} \mathbf{f}^{\sigma}$. Therefore the assertion must be formulated in terms of $N(\mathfrak{a})^{1 / 2} \lambda(\mathfrak{a})$.
21.5. Theorem. Let $\mathbf{f}$ be an eigenform contained in $\mathcal{S}_{k}(C)$ with $C$ as above; suppose that $\mathfrak{e}=\mathbf{c}$ and $\chi_{\mathbf{a}}(x)=x_{\mathbf{a}}^{[k]-e}\left|x_{\mathbf{a}}\right|^{i \kappa-[k]+e}$ with $e \in \mathbf{Z}^{\mathbf{a}}$ and $\kappa \in \mathbf{R}^{\mathbf{a}}$ such that $0 \leq e_{v} \leq 1$ for every $v \in \mathbf{a}$ and $\sum_{v \in \mathbf{a}} \kappa_{v}=0$. Put

$$
\mathcal{R}(s, \mathbf{f}, \chi)=\mathcal{Z}(s, \mathbf{f}, \chi) \prod_{v \in \mathbf{a}} \Gamma_{v}^{k, e}\left(\left(s+i \kappa_{v}\right) / 2\right)
$$

with $\Gamma_{v}^{k . e}$ defined as follows:

$$
\begin{gathered}
\Gamma_{v}^{k, e}(s)=\Gamma_{n}^{1}\left(s+\left(k_{v}+e_{v}-n-1\right) / 2\right) g\left(s, k_{v}-e_{v}\right) \\
g(s, a)= \begin{cases}\Gamma_{n}^{1}\left(s+\frac{a-n}{2}\right) & \text { if } a>n \\
\Gamma_{2 a+1-n}^{1}\left(s+\frac{a-n}{2}\right) \prod_{i=[a]+1}^{n-1} \Gamma\left(2 s-a-\frac{1}{2}\right) & \text { if }(n-2) / 2<a \leq n \\
\Gamma\left(s-\frac{n-2}{4}\right) \prod_{i=(n+1) / 2}^{n-1} \Gamma\left(2 s-a-\frac{1}{2}\right) \quad \text { if } a=(n-2) / 2 .\end{cases}
\end{gathered}
$$

Then $\mathcal{R}(s, \mathbf{f}, \chi)$ can be continued meromorphically to the whole s-plane with finitely many poles. Moreover each pole is simple. In particular, $\mathcal{R}$ is an entire function of $s$ if $\chi^{2} \neq 1$. If $\chi^{2}=1$, the poles are determined as follows: Let $\ell=\operatorname{Max}_{v \in \mathbf{a}}\left\{k_{v}-\right.$ $\left.e_{v}\right\}$. If $\ell>n, \mathcal{R}$ has no pole. If $\ell<n, \mathcal{R}$ has possible poles only in the set $\{(2 j+1) / 2 \mid j \in \mathbf{Z}, n+1 \leq j \leq 2 n+(1 / 2)-\ell\}$.

## 22. The largest possible pole of $\mathcal{Z}(s, \mathbf{f}, \chi)$

22.1. Put $\lambda_{n}=(n+1) / 2$ in Case SP and $\lambda_{n}=n$ in Case UT. We consider the set $S$ of (16.1a) and put

$$
\begin{align*}
S^{+} & =\left\{\xi \in S \mid \xi_{v}>0 \quad \text { for every } \quad v \in \mathbf{a}\right\}  \tag{22.1a}\\
S_{\mathbf{a}}^{+} & =\prod_{v \in \mathbf{a}} S_{v}^{+}, \quad S_{v}^{+}=\left\{\xi \in S_{v} \mid \xi>0\right\}
\end{align*}
$$

Given a subgroup $U$ of $G L_{n}(K)$, we define an equivalence relation $\sim \operatorname{in} S^{+}$with respect to $U$ by: $\sigma \sim \sigma^{\prime}$ if and only if $\sigma^{\prime}=a^{*} \sigma a$ with $a \in U$. We then denote by $S^{+} / U$ the set of all equivalence classes in this sense.

Take $\mathfrak{b}, \mathfrak{c}, \mathfrak{e}$, and $C$ as in $\S 19.1$; define $E^{\prime}$ by (19.2b). For $q \in G L_{n}(K)_{\mathbf{h}}$ and $\sigma \in S^{+}$put

$$
\begin{align*}
& U_{q}=G L_{n}(K) \cap q E^{*} q^{-1}, \quad E^{*}=G L_{n}(K)_{\mathbf{a}} E^{\prime},  \tag{22.2a}\\
& \nu_{\sigma, q}=\left[U_{\sigma, q}: 1\right]^{-1}, \quad U_{\sigma, q}=\left\{a \in U_{q} \mid a^{*} \sigma a=\sigma\right\} \tag{22.2~b}
\end{align*}
$$

Let $\mathbf{f} \in \mathcal{S}_{k}(C, \psi)$ and $\mathbf{g} \in \mathcal{M}_{l}\left(C^{\prime}, \psi^{\prime}\right)$ with integral or half-integral $k$ and $l$. Here $C^{\prime}$ is a group of the same type as $C$ with possibly different $\mathfrak{b}, \mathfrak{c}$, and $\mathfrak{e} ; \psi$ and $\psi^{\prime}$ are Hecke characters of $K$. We consider the Fourier coefficients $c_{\mathbf{f}}$ and $c_{\mathbf{g}}$ in the sense of Proposition 20.2. We assume that

$$
\begin{align*}
& \psi_{\mathbf{h}}^{\prime}(\operatorname{det}(a)) c_{\mathbf{g}}(\sigma, q a)=c_{\mathbf{g}}(\sigma, q)  \tag{22.3a}\\
& \left(\psi / \psi^{\prime}\right)_{\mathbf{a}}(x)=x_{\mathbf{a}}^{-t}\left|x_{\mathbf{a}}\right|^{t-i \kappa} \text { with every } a \in E^{\prime},  \tag{22.3b}\\
& t=\left\{\begin{array}{lc}
{[l]-[k]} & (\text { Rase SP }), \\
\left(l_{v}-l_{v \rho}-k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}} & \text { (Case UT). }
\end{array}\right. \tag{22.3c}
\end{align*}
$$

Notice that (20.9d) impies (22.3a) with $\mathbf{f}$ and $\psi$ in place of $\mathbf{g}$ and $\psi^{\prime}$. We then define, for $s \in \mathbf{C}$,

$$
\begin{equation*}
D_{q . \kappa}(s ; \mathbf{f}, \mathbf{g})=\sum_{\sigma \in S^{+} / U_{q}} \nu_{\sigma . q} c_{\mathbf{f}}(\sigma, q) \overline{c_{\mathbf{g}}(\sigma, q)} \operatorname{det}(\sigma)^{-s \mathbf{a}-h}, \text { where } \tag{22.4}
\end{equation*}
$$

$$
h=\left\{\begin{array}{lr}
(1 / 2)(k+l-i \kappa) & \text { (Case SP) }  \tag{22.4a}\\
(1 / 2)\left(k_{v}+k_{v \rho}+l_{v}+l_{v \rho}-i \kappa_{v}\right)_{v \in \mathbf{a}} & \text { (Case UT) }
\end{array}\right.
$$

This is well-defined. Indeed, if $a \in U_{q}$, then $\operatorname{det}(a) \in \mathfrak{r}^{\times}$and $c_{\mathbf{f}}\left(a^{*} \tau a, q\right)=\operatorname{det}(a)^{[k] \rho}$ $\cdot c_{\mathbf{f}}(\tau, a q)$ by $(20.9 \mathrm{~b}, \mathrm{c})$. Since $a q=q q^{-1} a q$ and $q^{-1} a q \in E^{*}$, by $(20.9 \mathrm{~b}, \mathrm{~d})$ we have

$$
\begin{equation*}
c_{\mathbf{f}}\left(a^{*} \tau a, q\right)=\psi_{\mathbf{a}}(\operatorname{det}(a)) \operatorname{det}(a)^{[k] \rho} c_{\mathbf{f}}(\tau, q) \text { for every } a \in U_{q} . \tag{22.5}
\end{equation*}
$$

A similar relation holds for $c_{\mathbf{g}}$. Combining these, we easily see that each term of (22.4) depends only on the equivalence class of $\sigma$ under $U_{q}$. Now the right-hand side of (22.4) is convergent for sufficiently large $\operatorname{Re}(s)$. This will be proven in §A6.7.
22.2. Proposition. Define $m, m^{\prime} \in \mathbf{Z}^{\mathbf{a}}$ as follows: $m=k$ and $m^{\prime}=l$ in Case $S P ; m=\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$ and $m^{\prime}=\left(l_{v}+l_{v \rho}\right)_{v \in \mathbf{a}}$ in Case UT. Then the following assertions hold:
(1) $D_{q, \kappa}(s ; \mathbf{f}, \mathbf{g})$ can be continued to a meromorphic function on the whole plane, which is holomorphic for $\operatorname{Re}(s)>0$. Moreover it is holomorphic at $s=0$ if $m \neq m^{\prime}$ or $\kappa \neq 0$.
(2) The right-hand side of (22.4) is absolutely convergent for $\operatorname{Re}(s)>0$ if $\mathbf{g}$ is a cusp form.
(3) Let $p=\operatorname{diag}[q, \widehat{q}]$; let $f_{p}$ and $g_{p}$ be the $p$-components of $\mathbf{f}$ and $\mathbf{g}$, respectively. If $m=m^{\prime}$ and $\kappa=0$, then $D_{q, 0}(s ; \mathbf{f}, \mathbf{g})$ has at most a simple pole at $s=0$ whose residue is a positive number times $\left\langle g_{p}, f_{p}\right\rangle$.

The proof will be completed in $\S 22.4$.
22.3. Let $f_{p}$ and $g_{p}$ be as in (3) above; take $\beta=1$ and $r=q$ in (20.9e). Then

$$
f_{p}(z)=\sum_{\sigma \in S_{+}} c_{\mathbf{f}}(\sigma, q) \mathbf{e}_{\mathbf{a}}^{n}(\sigma z), \quad g_{p}(z)=\sum_{\sigma \in S_{+}} c_{\mathbf{g}}(\sigma, q) \mathbf{e}_{\mathbf{a}}^{n}(\sigma z) .
$$

Take a congruence subgroup $\Gamma$ of $G_{1}$ so that $f_{p} \in \mathcal{M}_{k}(\Gamma)$ and $g_{p} \in \mathcal{M}_{l}(\Gamma)$; take also an r-lattice $L$ in $S$ so that $\left[\begin{array}{ll}1 & b \\ 0 & 1\end{array}\right] \in \Gamma$ for every $b \in L$. Let $X=S_{\mathrm{a}} / L$ and $Y=S_{\mathbf{a}}^{+} / U_{q}$, the latter being defined modulo the map $y \mapsto a \sigma a^{*}$ for $a \in U_{q}$. Write the variable $z$ on $\mathcal{H}^{n}$ in the form $z=x+i y$ with $x \in S_{\mathrm{a}}$ and $y \in S_{\mathbf{a}}^{+}$; define a measure $d x$ on $S_{\mathbf{a}}$ by $d x=\prod_{v \in \mathbf{a}} d x_{v}$ with $d x_{v}$ defined on $S_{v}$ as in $\S 5.12$. Clearly

$$
\begin{equation*}
\int_{X} f_{p}(x+i y) \overline{g_{p}(x+i y)} d x=\operatorname{vol}(X) \sum_{\sigma \in S^{+}} c_{\mathbf{f}}(\sigma, q) \overline{c_{\mathbf{g}}(\sigma, q)} \mathbf{e}_{\mathbf{a}}^{n}(2 i \sigma y) \tag{22.6}
\end{equation*}
$$

We recall a well-known formula

$$
\begin{equation*}
\int_{S_{v}^{+}} e^{-\operatorname{tr}(y \tau)} \operatorname{det}(y)^{s-\lambda_{n}} d y=\Gamma_{n}^{\iota}(s) \operatorname{det}(\tau)^{-s} \quad\left(\operatorname{Re}(s)>\lambda_{n}-1 ; \tau \in S_{v}^{+}\right) \tag{22.7}
\end{equation*}
$$

where $\Gamma_{n}^{\iota}$ is defined by (16.47). Put

$$
H(y)=\sum_{\sigma \in \mathcal{T}} \nu_{\sigma, q} M(\sigma, y), \quad M(\sigma, y)=c_{\mathbf{f}}(\sigma, q) \overline{c_{\mathbf{g}}(\sigma, q)} \mathbf{e}_{\mathbf{a}}^{n}(2 i \sigma y) \operatorname{det}(y)^{s \mathbf{a}+h}
$$

where $\mathcal{T}$ is a fixed complete set of representatives for $S^{+} / U_{q}$. Since $M\left(a^{*} \sigma a, y\right)=$ $M\left(\sigma, a y a^{*}\right)$ for every $a \in U_{q}$ by (22.5), putting $d^{\prime} y=\operatorname{det}(y)^{-\lambda_{n} \mathbf{a}} d y$, we have

$$
\int_{Y} \sum_{\sigma \in S^{+}} M(\sigma, y) d^{\prime} y=\int_{Y} \sum_{a \in U_{q}} \sum_{\sigma \in \mathcal{T}} \nu_{\sigma, q} M\left(a^{*} \sigma a, y\right) d^{\prime} y=\int_{Y} \sum_{a \in U_{q}} H\left(a y a^{*}\right) d^{\prime} y
$$

$$
\begin{gathered}
=\nu \int_{S_{\mathbf{a}}^{+}} H(y) d^{\prime} y=\nu \sum_{\sigma \in \mathcal{T}} \nu_{\sigma, q} \int_{S_{\mathbf{a}}^{+}} M(\sigma, y) d^{\prime} y=\nu \cdot \Gamma((s)) D_{q, \kappa}(s ; \mathbf{f}, \mathbf{g}) \\
\text { where } \Gamma((s))=\prod_{v \in \mathbf{a}} \Gamma_{n}^{\iota}\left(s+h_{v}\right)(4 \pi)^{-n\left(s+h_{v}\right)}
\end{gathered}
$$

and $\nu$ is the number of the roots of unity $\zeta$ in $K$ such that $\zeta-1 \in \mathfrak{e}$; we employed (22.7) in the last step. Termwise integration can be justified for sufficiently large $\operatorname{Re}(s)$ because of the convergence of (22.4). Combining this with (22.6), we obtain

$$
\nu \cdot \operatorname{vol}(X) \Gamma((s)) D_{q, \kappa}(s ; \mathbf{f}, \mathbf{g})=\int_{X} \int_{Y} f_{p}(x+i y) \overline{g_{p}(x+i y)} \delta(z)^{h+s \mathbf{a}-\lambda_{n} \mathbf{a}} d x d y
$$

for sufficiently large $\operatorname{Re}(s)$, where $\delta(x+i y)=\left(\operatorname{det}\left(y_{v}\right)\right)_{v \in \mathbf{a}}$. Put $\Gamma^{P}=\Gamma \cap P$ as in (17.2). Then $\left[\begin{array}{ll}a & b \\ 0 & \widehat{a}\end{array}\right] \in \Gamma^{P}$ for every $b \in L$ and every $a$ in a suitable subgroup $U^{\prime}$ of $U_{q}$ of finite index. Therefore $X \times Y$ is "commensurable with" $\Gamma^{P} \backslash \mathcal{H}$. Now we have $\delta^{-2 \lambda_{n} \mathbf{a}} d x d y=c \mathbf{d} z=c \prod_{v \in \mathbf{a}} \mathbf{d} z_{v}$ with $\mathbf{d} z_{v}$ as in Lemma 3.4; $c=1$ in Case SP and $c=2^{n(1-n)[F: \mathbf{Q}]}$ in Case UT as explained in §5.12. Thus the last double integral is a positive rational number times

$$
\begin{equation*}
\int_{Z} f_{p}(z) \overline{g_{p}(z)} \delta(z)^{h+s \mathbf{a}+\lambda_{n} \mathbf{a}} \mathbf{d} z \quad\left(Z=\Gamma^{P} \backslash \mathcal{H}^{n}\right) \tag{22.8}
\end{equation*}
$$

Clearly $j_{\gamma}^{k-l} j_{\gamma}^{l}=j_{\gamma}^{k}$ if both $k$ and $l$ are integral, but the case of half-integral weights is not so simple. However, by (16.17) and Theorem 6.8 (5), the equality always holds for $\gamma$ in a suitable congruence subgroup $\Gamma$ of $G_{1}$. With this choice of $\Gamma$, we have $j_{\gamma}^{k} \overline{j_{\gamma}^{l}}=\left(\overline{j_{\gamma}^{k-l}}\right)^{-1}\left|j_{\gamma}^{k}\right|^{2}$ for $\gamma \in \Gamma$. Thus, putting $\mathcal{F}(z)=f_{p}(z) \overline{g_{p}(z)} \delta^{h+s^{\prime} \text { a }}$ with $s^{\prime}=s+\lambda_{n}$ and defining $m$ and $m^{\prime}$ as in our proposition, we find that

$$
\mathcal{F} \circ \gamma=\mathcal{F} \cdot\left|j_{\gamma}\right|^{2 m-2 h-2 s^{\prime} \mathbf{a}} \overline{\left(j_{\gamma}^{k-l}\right)}-1 \text { for every } \gamma \in \Gamma
$$

Let $\mathfrak{D}=\Gamma \backslash \mathcal{H}^{n}$ and $\mathcal{R}=\Gamma^{P} \backslash \Gamma$. Choosing $\Gamma$ suitably, we may assume that $\Gamma \cap K^{\times}=$ $\{1\}$. Then (22.8) equals

$$
\int_{Z} \mathcal{F}(z) \mathbf{d} z=\int_{\mathfrak{D}} \sum_{\gamma \in \mathcal{R}} \mathcal{F}(\gamma(z)) \mathbf{d} z=\int_{\mathfrak{D}} \mathcal{F}(z)\left\{\sum_{\gamma \in \mathcal{R}}\left|j_{\gamma}\right|^{2 m-2 h-2 s^{\prime} \mathbf{a}} \overline{\left(j_{\gamma}^{k-l}\right)}-1\right\} \mathbf{d} z .
$$

The last sum over $\mathcal{R}$ can be written $\delta^{m-h-s^{\prime}} \mathbf{a} \overline{E\left(z, \bar{s}+\lambda_{n} ; m-m^{\prime}, \kappa, \Gamma\right)}$ with $E(\cdots)$ of (17.3). Thus

$$
\begin{equation*}
\Gamma((s)) D_{q, \kappa}(s ; \mathbf{f}, \mathbf{g})=A \int_{\mathfrak{D}} f_{p}(z) \overline{g_{p}(z)} \overline{E\left(z, \bar{s}+\lambda_{n} ; m-m^{\prime}, \kappa, \Gamma\right)} \delta(z)^{m} \mathrm{~d} z \tag{22.9}
\end{equation*}
$$

with a constant $A$, which is a positive rational number times $\operatorname{vol}(X)^{-1}$.
22.4. Proof of Proposition 22.2. By Lemma 17.2 (4), $E\left(z, s ; m-m^{\prime}, \kappa, \Gamma\right)$ is holomorphic for $\operatorname{Re}(s)>\lambda_{n}$. Moreover, by a well-known principle, if it has a pole of order $t$ at $s_{0}$, then $\left(s-s_{0}\right)^{t} E$ as a function of $z$ is slowly increasing at every cusp. If $\mathbf{f} \neq 0$, (6.42) shows that $m_{v}>\lambda_{n}-1$ for every $v \in \mathbf{a}$, and hence $\prod_{v \in \mathbf{a}} \Gamma_{n}^{\iota}\left(m_{v}\right)^{-1} \neq 0$. Therefore (1) and (3) follow immediately from (22.9) combined with Lemma 17.2 (4). To prove (2), we observe that if the series for $D_{q, 0}(s ; \mathbf{f}, \mathbf{f})$ and $D_{q, 0}(s ; \mathbf{g}, \mathbf{g})$ for real $s$ are convergent, then the Cauchy-Schwarz inequality gives

$$
\sum_{\sigma \in S^{+} / U_{q}} \nu_{\sigma, q}\left|c_{\mathbf{f}}(\sigma, q) c_{\mathbf{g}}(\sigma, q) \operatorname{det}(\sigma)^{-s \mathbf{a}-h}\right| \leq\left[D_{q, 0}(s ; \mathbf{f}, \mathbf{f}) D_{q .0}(s ; \mathbf{g}, \mathbf{g})\right]^{1 / 2}
$$

and therefore assertion (2) in the general case follows from the special case $\mathbf{f}=$ g. Now we know that $D_{q, 0}(s ; \mathbf{f}, \mathbf{f})$ is holomorphic for $\operatorname{Re}(s)>0$. Since all its
coefficients are nonnegative, the series defining $D_{q, 0}(s ; \mathbf{f}, \mathbf{f})$ must be convergent for $\operatorname{Re}(s)>0$. This completes the proof.
22.5. Now given $\mathbf{f} \in \mathcal{S}_{k}(C, \psi)$, a Hecke character $\chi$ of $K, \tau \in S^{+}$, and $r \in$ $G L_{n}(K)_{\mathbf{h}}$, we define a formal Dirichlet series $\mathcal{D}(\mathbf{f}, \chi)$ and an ordinary Dirichlet series $D(s, \mathbf{f}, \chi)$ by

$$
\begin{equation*}
\mathcal{D}(\mathbf{f}, \chi)=\sum_{x \in B^{\prime} / E^{\prime}} \psi(\operatorname{det}(r x)) \chi^{*}(\operatorname{det}(x) \mathfrak{r}) c_{\mathbf{f}}(\tau, r x)\left|\operatorname{det}\left(x x^{*}\right)\right|_{F}^{-\lambda_{n}}[\operatorname{det}(x) \mathfrak{r}] \tag{22.10}
\end{equation*}
$$

$$
\begin{equation*}
D(s, \mathbf{f}, \chi)=\sum_{x \in B^{\prime} / E^{\prime}} \psi(\operatorname{det}(r x)) \chi^{*}(\operatorname{det}(x) \mathfrak{r}) c_{\mathbf{f}}(\tau, r x)\left|\operatorname{det}\left(x x^{*}\right)\right|_{F}^{-\lambda_{n}}|\operatorname{det}(x)|_{K}^{s} \tag{22.11}
\end{equation*}
$$

where $B^{\prime}$ and $E^{\prime}$ are as in (19.2b), and $\left|\left.\right|_{F}\right.$ is the idele norm in $F_{\mathbf{A}}^{\times}$. These depend on $r$ and $\tau$, but we fix them in the following treatment. Let

$$
\begin{equation*}
G L_{n}(K)_{\mathbf{A}}=\bigsqcup_{q \in Q} G L_{n}(K) q E^{*} \tag{22.12}
\end{equation*}
$$

with a finite subset $Q$ of $G L_{n}(K)_{\mathbf{h}}$; put then

$$
\begin{gathered}
X_{q}=G L_{n}(K) \cap r B^{*} q^{-1}, \quad X_{\sigma, q}=\left\{\xi \in X_{q} \mid \sigma=\xi^{*} \tau \xi\right\} \quad\left(q \in Q, \sigma \in S^{+}\right), \\
B^{*}=B^{\prime} \cdot G L_{n}(K)_{\mathbf{a}} .
\end{gathered}
$$

Given $x \in B^{*}$, take $\xi \in G L_{n}(K)$ and $q \in Q$ so that $r x \in \xi q E^{*}$. Then $x \in r^{-1} \xi q E^{*}$ and $\xi \in X_{q}$. From this we easily see that $B^{\prime} / E^{\prime}$ can be given by $\bigsqcup_{q \in Q}\left\{r^{-1} \xi_{\mathbf{h}} q \mid\right.$ $\left.\xi \in X_{q} / U_{q}\right\}$. Therefore we have

$$
\begin{align*}
& \quad \mathcal{D}(\mathbf{f}, \chi)=\sum_{q \in Q}\left|\operatorname{det}\left(r^{-1} \widehat{r} q q^{*}\right)\right|_{F}^{-\lambda_{n}}  \tag{22.13}\\
& \sum_{\xi \in X_{q} / U_{q}} c_{\mathbf{f}}(\tau, \xi q) \psi_{\mathbf{h}}(\operatorname{det}(\xi q)) \chi^{*}\left(\operatorname{det}\left(r^{-1} \xi q\right) \mathfrak{r}\right)|\operatorname{det}(\xi)|^{2 \lambda_{n} \mathbf{a}}\left[\operatorname{det}\left(r^{-1} \xi q\right) \mathfrak{r}\right]
\end{align*}
$$

22.6. Take $\mu \in \mathbf{Z}^{\mathbf{b}}$ under the following conditions: $\mu_{v} \geq 0$ for every $v \in \mathbf{b}$; $\mu_{v} \leq 1$ for every $v \in \mathbf{a}$ in Case SP; $\mu_{v} \mu_{v \rho}=0$ for every $v \in \mathbf{a}$ in Case UT. Given $\chi$ as above, let $\mathfrak{f}$ be the conductor of $\chi$. We define a weight $l$ and a Hecke character $\psi^{\prime}$ of $K$ by

$$
\begin{align*}
& l= \begin{cases}\mu+(n / 2) \mathbf{a} & (\text { Case SP }) \\
\mu+n \mathbf{a}\end{cases}  \tag{22.14a}\\
& \psi^{\prime}= \begin{cases}\chi^{-1} \rho_{\tau} & (\text { Case UT }) \\
\chi^{-1} \varphi^{-n} & (\text { Case SP })\end{cases}  \tag{22.14b}\\
& \text { (Case UT) }
\end{align*}
$$

Here $\rho_{\tau}$ is the Hecke character of $F$ corresponding to the extension $F\left(c^{1 / 2}\right) / F$ with $c=(-1)^{[n / 2]} \operatorname{det}(2 \tau) ; \varphi$ is a Hecke character of $K$ such that $\varphi(y)=y^{-\mathbf{a}}|y|^{\mathbf{a}}$ for $y \in K_{\mathbf{a}}^{\times}$and the restriction of $\varphi$ to $F_{\mathbf{A}}^{\times}$is the Hecke character of $F$ corresponding to $K / F$. Such a $\varphi$ always exists, but not necessarily unique; see Lemma A5.1. We fix one such $\varphi$ in the following treatment.

As will be shown in $\S$ A 5.5 , there exists an element $\mathbf{g} \in \mathcal{M}_{l}\left(C^{\prime}, \psi^{\prime}\right)$ such that

$$
\begin{align*}
c_{\mathbf{g}}(\sigma, q)= & |\operatorname{det}(q)|_{K}^{n / 2} \psi^{\prime}(\operatorname{det}(q))^{-1}  \tag{22.15}\\
& \cdot \sum_{\xi \in X_{\sigma, q}} \bar{\chi}_{\mathbf{a}}(\operatorname{det}(\xi)) \bar{\chi}^{*}\left(\operatorname{det}\left(r^{-1} \xi q\right) \mathfrak{r}\right) \operatorname{det}(\xi)^{\mu \rho}
\end{align*}
$$

for $\sigma \in S^{+}$, where

$$
C^{\prime}=\left\{\alpha \in D\left[\mathfrak{b}_{1}^{-1} \mathfrak{e}, \mathfrak{b}_{1} \mathfrak{c}_{1}\right] \mid a_{\alpha}-1 \prec \mathfrak{r e}\right\}
$$

with the same $\mathfrak{e}$ as before and some $\mathfrak{b}_{1}$ and $\mathfrak{c}_{1} ; \mathfrak{c}_{1}$ is divisible by $\mathfrak{e}, \mathfrak{f}$, and the conductor of $\rho_{\tau}$ or $\varphi$.

We take this $\mathbf{g}$ to be that in (22.4); naturally (22.3b, c) must be satisfied, which is so only for a suitable choice of $\mu$. In fact, given $\mathbf{f} \in \mathcal{S}_{k}(C, \psi)$ and $\chi$, define $\psi^{\prime}$ as above, and put $(\psi \chi)_{\mathbf{a}}(x)=x_{\mathbf{a}}^{-t^{\prime}}\left|x_{\mathbf{a}}\right| t^{t^{\prime}-i \kappa}$ with $t^{\prime} \in \mathbf{Z}^{\mathbf{a}}$ and $\kappa \in \mathbf{R}^{\mathbf{a}}$. Then $\left(\psi / \psi^{\prime}\right)_{\mathbf{a}}(x)=x_{\mathbf{a}}^{-t}\left|x_{\mathbf{a}}\right|^{t-i \kappa}$ with $t=t^{\prime}+[n / 2] \mathbf{a}$ in Case SP and $t=t^{\prime}+n \mathbf{a}$ in Case UT. Now take $\mu \in \mathbf{Z}^{\mathbf{b}}$ as follows: In Case SP define $\mu_{v}$ by the conditions $0 \leq \mu_{v} \leq 1$ and $\mu-[k]-t^{\prime} \in 2 \mathbf{Z}^{\text {a }}$. In Case UT, put

$$
\begin{align*}
& \mu_{v}=t_{v}^{\prime}-k_{v \rho}+k_{v} \quad \text { and } \quad \mu_{v \rho}=0 \quad \text { if } \quad t_{v}^{\prime} \geq k_{v \rho}-k_{v}  \tag{22.15a}\\
& \mu_{v}=0 \quad \text { and } \quad \mu_{v \rho}=k_{v \rho}-k_{v}-t_{v}^{\prime} \quad \text { if } \quad t_{v}^{\prime}<k_{v \rho}-k_{v} . \tag{22.15b}
\end{align*}
$$

We consider $\mathbf{g}$ with this $\mu$. Then (22.3b, c) are satisfied.
Let $\mathcal{D}_{q}(\mathbf{f}, \mathbf{g})$ denote the formal Dirichlet series obtained from the right-hand side of (22.4) with $\operatorname{det}(\sigma)^{-h}[\operatorname{det}(\sigma) \mathfrak{r}]$ in place of $\operatorname{det}(\sigma)^{-s \mathbf{a}-h}$. Then

$$
\begin{align*}
& \operatorname{det}(\tau)^{h}[\operatorname{det}(\tau) \mathfrak{r}]^{-1}|\operatorname{det}(q)|_{K}^{-n / 2} \psi^{\prime}\left(\operatorname{det}(q)^{-1}\right) \mathcal{D}_{q}(\mathbf{f}, \mathbf{g})  \tag{22.16}\\
& =\sum_{\xi \in X_{q} / U_{q}} \psi_{\mathbf{h}}(\operatorname{det}(\xi)) \chi^{*}\left(\operatorname{det}\left(r^{-1} \xi q\right) \mathfrak{r}\right) c_{\mathbf{f}}(\tau, \xi q)|\operatorname{det}(\xi)|^{-(n / u) \mathbf{a}}\left[\operatorname{det}\left(\xi \xi^{*}\right) \mathfrak{r}\right]
\end{align*}
$$

where $u=2 /[K: F]$. To see this, take complete sets of representatives $R$ and $Y_{\sigma, q}$ for $S^{+} / U_{q}$ and $X_{\sigma, q} / U_{\sigma, q}$, respectively. Then we easily see that $\bigsqcup_{\sigma \in R} Y_{\sigma, q}$ gives $X_{q} / U_{q}$. In Case UT, for $x \in K^{\times}$we can easily verify that

$$
\begin{equation*}
(\psi \chi)_{\mathbf{a}}(x)^{-1} x^{-k \rho}=x^{t^{\prime}-k \rho}|x|^{i \kappa-t^{\prime}}=x^{\mu}|x|^{n \mathbf{a}-2 h} . \tag{22.17}
\end{equation*}
$$

Here we have to remember that $k \rho$ and $h$ belong to $\mathbf{Z}^{\mathbf{b}} ; t^{\prime}$ and $h$ ar elements of $\mathbf{Z}^{\mathbf{a}}$ viewed also as elements of $\mathbf{Z}^{\mathbf{b}}$. To make our formulas short, let us put ${ }^{d \xi}=\operatorname{det}(\xi)$ temporarily. Then the sum over $X_{q} / U_{q}$ in (22.16) in Case UT equals $\sum_{\sigma \in R} \nu_{\sigma, q} A_{\sigma}$ with

$$
\begin{aligned}
A_{\sigma} & =\left.\left.\sum_{\xi \in X_{\sigma, q}} \psi_{\mathbf{a}}\left({ }^{d} \xi\right)^{-1} c_{\mathbf{f}}(\tau, \xi q)\right|^{d} \xi\right|^{-n \mathbf{a}} \chi^{*}\left({ }^{d}\left(r^{-1} \xi q\right) \mathfrak{r}\right)\left[^{d}\left(\xi \xi^{*}\right) \mathfrak{r}\right] \\
& \left.=\left.\left.\sum_{\xi} \psi_{\mathbf{a}}\left({ }^{d} \xi\right)^{-1} c_{\mathbf{f}}(\sigma, q)\left({ }^{(d} \xi\right)^{-k \rho}\right|^{d} \xi\right|^{-n \mathbf{a}} \chi^{*}\left({ }^{d}\left(r^{-1} \xi q\right) \mathfrak{r}\right){ }^{d}\left(\xi \xi^{*}\right) \mathfrak{r}\right] \\
& =\left.\left.c_{\mathbf{f}}(\sigma, q) \sum_{\xi}(\psi \chi) \mathbf{a}\left({ }^{( } \xi\right)^{-1}\left(\frac{d}{}\right)^{-k \rho}\right|^{d} \xi\right|^{-n \mathbf{a}} \chi_{\mathbf{a}}\left({ }^{(d)} \xi\right) \chi^{*}\left({ }^{d}\left(r^{-1} \xi q\right) \mathfrak{r}\right)\left[^{d}\left(\xi \xi^{*}\right) \mathfrak{r}\right] \\
& =\left.\left.\right|^{d}\left(\tau^{-1} \sigma\right)\right|^{-h}\left[{ }^{d}\left(\tau^{-1} \sigma\right) \mathfrak{r}\right] c_{\mathbf{f}}(\sigma, q) \sum_{\xi} \chi_{\mathbf{a}}\left({ }^{d} \xi\right) \chi^{*}\left({ }^{d}\left(r^{-1} \xi q\right) \mathfrak{r}\right)\left({ }^{d} \xi\right)^{\mu},
\end{aligned}
$$

where we employed (22.17) in the last step. By (22.15) the last sum over $\xi$ equals $|\operatorname{det}(q)|_{K}^{-n / 2} \psi^{\prime}(\operatorname{det}(q))^{-1} \overline{c_{\mathbf{g}}(\sigma, q)}$. Therefore we obtain (22.16) in Case UT. Case SP can be handled in a similar and simpler way.

Now (22.16) together with (22.13) gives

$$
\begin{align*}
\text { a) } & \mathcal{D}\left(N(\mathfrak{a})^{-s_{0}}\left[\mathfrak{a} \mathfrak{a}^{\rho}\right]\right)  \tag{22.18a}\\
= & \operatorname{det}(\tau)^{h}|\operatorname{det}(r)|_{K}^{-n / 2} \sum_{q \in Q}\left(\psi / \psi^{\prime}\right)(\operatorname{det}(q))\left[\operatorname{det}\left(r^{-1} \widehat{r} q q^{*} \tau^{-1}\right) \mathfrak{r}\right] \mathcal{D}_{q}(\mathbf{f}, \mathbf{g}), \\
\text { b) } \quad & D\left(u s+s_{0}, \mathbf{f}, \chi\right)  \tag{22.18b}\\
= & \operatorname{det}(\tau)^{s \mathbf{a}+h}|\operatorname{det}(r)|_{K}^{-u s-n / 2} \sum_{q \in Q}\left(\psi / \psi^{\prime}\right)(\operatorname{det}(q))\left|\operatorname{det}\left(q q^{*}\right)\right|_{F}^{s} D_{q . \kappa}(s ; \mathbf{f}, \mathbf{g}),
\end{align*}
$$

where $u=2 /[K: F], s_{0}=(3 n / 2)+u-1$, and the left-hand side of (22.18a) means the formal series obtained from $\mathcal{D}(\mathbf{f}, \chi)$ by substituting $N(\mathfrak{a})^{-s_{0}}\left[\mathfrak{a} \mathfrak{a}^{\rho}\right]$ for [a].
22.7. Lemma. Given a formal Dirichlet series $\sum_{\mathfrak{a} \subset \mathfrak{r}} c_{\mathfrak{a}}[\mathfrak{a}]$ with $c_{\mathfrak{a}} \in \mathbf{C}$, suppose that $\sum_{\mathfrak{a} \subset \mathfrak{r}}\left|c_{\mathfrak{a}} N(\mathfrak{a})^{-s}\right|<\infty$ for $\operatorname{Re}(s)>\alpha$ and also that it can be decomposed formally into an Euler product in the sense that $\sum_{\mathfrak{a}} c_{\mathfrak{a}}[\mathfrak{a}]=\prod_{\mathfrak{p}} V_{\mathfrak{p}}([\mathfrak{p}])^{-1}$ with complex polynomials $V_{\mathfrak{p}}(x)$ such that $V_{\mathfrak{p}}(0)=1$ defined for all prime ideals $\mathfrak{p}$ in $K$. Then the infinite product $\prod_{\mathfrak{p}} V_{\mathfrak{p}}\left(\varphi(\mathfrak{p}) N(\mathfrak{p})^{-s}\right)^{-1}$, for any $\mathbf{T}$-valued ideal character $\varphi$, is convergent to the Dirichlet series $\sum_{\mathfrak{a} \subset \mathfrak{r}} c_{\mathfrak{a}} \varphi(\mathfrak{a}) N(\mathfrak{a})^{-s}$ and nonvanishing for $\operatorname{Re}(s)>\alpha$.

Proof. We can formally put $V_{\mathfrak{p}}(x)^{-1}=1+\sum_{n=1}^{\infty} b_{\mathfrak{p}, n} x^{n}$ with $b_{\mathfrak{p}, n} \in \mathbf{C}$. Then $\sum_{\mathfrak{p}} \sum_{n=1}^{\infty} b_{\mathfrak{p}, n} N(\mathfrak{p})^{-n s}$ is a partial series of $\sum_{\mathfrak{a}} c_{\mathfrak{a}} N(\mathfrak{a})^{-s}$ and therefore $\sum_{\mathfrak{p}} \sum_{n=1}^{\infty}\left|b_{\mathfrak{p}, n} \varphi(\mathfrak{p})^{n} N(\mathfrak{p})^{-n s}\right|<\infty$ for $\operatorname{Re}(s)>\alpha$. Thus the infinite product must be convergent. Clearly each factor $V_{\mathfrak{p}}\left(\varphi(\mathfrak{p}) N(\mathfrak{p})^{-s}\right)^{-1}$ is nonvanishing, and hence we obtain our lemma.
22.8. Lemma. Given $(\psi, \tau)$ as above and an integral ideal $\mathfrak{x}$ in $F$, there exists a Hecke character $\chi$ and $\mu \in \mathbf{Z}^{\mathbf{b}}$ with the following properties: (i) our discussion of $\S 22.6$ is applicable to $(\chi, \mu)$; (ii) $\mu \neq 0$; (iii) the conductor of $\chi$ is prime to $\mathfrak{x}$.

Proof. In Case SP take $\mu \in \mathbf{Z}^{\mathbf{a}}, \neq 0$, so that $0 \leq \mu_{v} \leq 1$, and put $t^{\prime}=\mu-[k]$; in Case UT take $\mu \in \mathbf{Z}^{\mathbf{b}}, \neq 0$, so that $\mu_{v} \geq 0$ and $\mu_{v \rho}=0$ for every $v \in \mathbf{a}$, and put $t^{\prime}=\left(\mu_{v}-k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$. We can find an integral $\mathfrak{r}$-ideal $\mathfrak{a}$ prime to $\mathfrak{x}$ so that if $\zeta$ is a root of unity in $K$ and $\zeta-1 \in \mathfrak{a}$, then $\zeta=1$. Put $W=\{e \in$ $\left.\mathfrak{r}^{\times} \mid e-1 \in \mathfrak{a}\right\}$. Then the map $e \mapsto\left(\log \left|e_{v}\right|\right)_{v \in \mathbf{a}}$ gives an isomorphism of $W$ onto a $\mathbf{Z}$-lattice in $\left\{x \in \mathbf{R}^{\mathbf{a}} \mid \sum_{v \in \mathbf{a}} x_{v}=0\right\}$. Now $e \mapsto \psi_{\mathbf{a}}(e)^{-1}|e|^{t^{\prime}} e^{-t^{\prime}}$ defines a $\mathbf{T}$-valued character of $W$. Therefore we can easily find $\kappa \in \mathbf{R}^{\mathbf{a}}$ such that $\sum_{v \in \mathbf{a}} \kappa_{v}=0$ and $\psi_{\mathbf{a}}(e)^{-1}|e|^{t^{\prime}} e^{-t^{\prime}}=|e|^{i \kappa}$ for every $e \in W$. Now $\psi_{\mathbf{a}}(x)=x^{c}|x|^{-c-i \kappa^{\prime}}$ with $c \in \mathbf{Z}^{\mathbf{a}}$ and $\kappa^{\prime} \in \mathbf{R}^{\mathbf{a}}, \sum_{v \in \mathbf{a}} \kappa_{v}^{\prime}=0$. Therefore, by [S97, Lemma 11.14 (1)] there exists a Hecke character $\chi$ of $K$ whose conductor divides $\mathfrak{a}$ and such that $\chi_{\mathbf{a}}(x)=$ $\psi_{\mathbf{a}}(x)^{-1} x_{\mathbf{a}}^{-t^{\prime}}\left|x_{\mathbf{a}}\right| t^{t^{\prime}-i \kappa}$. This proves our lemma.
22.9. We are now ready to prove Theorem 20.13. Given an eigenform $\mathbf{f} \in$ $\mathcal{S}_{k}(C, \psi)$, take $\mu$ and $\chi$ as in Lemma 22.8 with an arbitrary choice of $\mathfrak{x}$. Since $\mathbf{f}$ is a cusp form, by Lemma 20.8 and Theorem 20.9 we can find $r \in G L_{n}(K)_{\mathbf{h}}$ and $\tau \in S^{+}$such that $c_{\mathbf{f}}(\tau, r) \neq 0$ and (20.19) holds. Then we consider (22.18a, b) for these $\tau, r, \chi$. Observe that $D(s, \mathbf{f}, \chi)$ can be obtained from $D(\tau, r ; \mathbf{f})$ of Theorem 20.4 by substituting $\chi^{*}(t \mathfrak{r})\left(\psi / \psi_{\mathrm{c}}\right)(t)|t|_{K}^{s}$ for $[t \mathbf{t r}], t \in K_{\mathrm{h}}^{\times}$, and multiplying by $\left(\psi / \psi_{\mathrm{c}}\right)(\operatorname{det}(r))$. Let $\mathcal{Z}^{\prime}(s)$ denote the Euler product obtained from $\mathcal{Z}(s, \mathbf{f}, \chi)$ of (20.21) by substituting $\left(\psi / \psi_{\mathfrak{c}}\right)\left(\pi_{\mathfrak{P}}\right) N(\mathfrak{P})^{-s}$ for $N(\mathfrak{P})^{-s}$ for each prime ideal $\mathfrak{P}$ in $K$, where $\pi_{\mathfrak{P}}$ is the prime element of $K_{\mathfrak{P}}$. Then (20.19) gives an equality between formal Dirichlet series which leads to

$$
\begin{equation*}
\psi_{\mathbf{c}}(\operatorname{det}(r)) c_{\mathbf{f}}(\tau, r) \mathcal{Z}^{\prime}(s)=D(s, \mathbf{f}, \chi) P(s) \Lambda(s) \tag{22.19}
\end{equation*}
$$

Here $P$ is a finite Dirichlet series with constant term 1 , and $\Lambda$ is the product of certain $L$-functions obtained from $\mathfrak{L}_{0} \prod_{v \nmid c} h_{v}([\mathfrak{p r}])^{-1}$ of (20.19) by substituting $\chi^{*}(\mathfrak{p r}) \psi^{*}(\mathfrak{p r}) N(\mathfrak{p r})^{-s}$ for $[\mathfrak{p r}]$ for each prime ideal $\mathfrak{p}$ in $F$.

Put $\mathcal{D}_{q}(\mathbf{f}, \mathbf{g})=\sum_{\sigma \in R} a_{\sigma}^{q}[\operatorname{det}(\sigma) \mathfrak{r}]$ and $\mathcal{D}(\mathbf{f}, \chi)=\sum_{\mathfrak{a}} b_{\mathfrak{a}}[\mathfrak{a}]$ with $a_{\sigma}^{q}, b_{\mathfrak{a}} \in \mathbf{C}$. Then $D(s, \mathbf{f}, \chi)=\sum_{\mathfrak{a}} b_{\mathfrak{a}} N(\mathfrak{a})^{-s}$, and (22.18a) shows that $\sum_{\mathfrak{a}} b_{\mathfrak{a}} N(\mathfrak{a})^{-s_{0}}\left[\mathfrak{a} a^{\rho}\right]=$ $\sum_{q \in Q} d_{q} \sum_{\sigma \in R} a_{\sigma}^{q}[\operatorname{det}(\sigma) \mathfrak{r}]$ with $d_{q} \in \mathbf{C}$. Our $\mathbf{g}$ will be given in $\S A 5.5$ as a certain theta series $\theta(z, \lambda)$. Since $\mu \neq 0, \mathbf{g}$ is a cusp form by (A3.16) or (A5.5). Therefore

Proposition 22.2 (2) applied to $D_{q, \kappa}(s ; \mathbf{f}, \mathbf{g})$ in (22.18b) shows that $\sum_{\mathfrak{a}}\left|b_{\mathfrak{a}} N(\mathfrak{a})^{-s}\right|$ $<\infty$ for $\operatorname{Re}(s)>s_{0}$. The same holds obviously for $P$, and also for $\Lambda$, as can easily be seen from the explicit forms of $\mathfrak{L}_{0}$ and $h_{v}([\mathfrak{p r}])$. Thus the Dirichlet series expressing $\mathcal{Z}^{\prime}(s)$ has the same property, and clearly the same is true for $\mathcal{Z}(s, \mathbf{f}, \chi)$. Now given an arbitrary Hecke character $\omega$ of $K$, Lemma 22.7 shows that the Euler product for $\mathcal{Z}(s, \mathbf{f}, \omega)$ is absolutely convergent and nonvanishing for $\operatorname{Re}(s)>s_{0}$, if we remove the Euler $\mathfrak{P}$-factors for $\mathfrak{P}$ dividing the conductor of $\chi$. By Lemma 22.8, we can choose $\chi$ so that the conductor of $\chi$ is prime to any integral ideal in $F$, which means that the Euler $\mathfrak{P}$-factor of $\mathcal{Z}(s, \mathbf{f}, \omega)$ for an arbitrary $\mathfrak{P}$ is nonvanishing for $\operatorname{Re}(s)>s_{0}$. Consequently the whole Euler product for $\mathcal{Z}(s, \mathbf{f}, \omega)$ is nonvanishing for $\operatorname{Re}(s)>s_{0}$. This proves the second part of Theorem 20.13.

To prove the first part concerning meromorphic continuation, observe that the above $\mathcal{Z}^{\prime}$ coincides with $\mathcal{Z}(s, \mathbf{f}, \chi \psi)$ if we remove, if necessary, finitely many Euler factors. Here we can take an arbitrary $\chi$; we do not have to assume $\mu \neq 0$; still (22.19) holds. Therefore meromorphic continuation of $\mathcal{Z}(s, \mathbf{f}, \chi)$ for an arbitrary $\chi$ follows from Proposition 22.2 (1), (22.18b), and (22.19).

In the above proof we employed formal Dirichlet series and stated equality (22.18a) in order to emphasize that (22.18b) is not merely an equality between two functions, but also the equalities between the formal Dirichlet series, which is essential for the desired convergence of the series.

Theorem 20.13 concerns integral $k$. As noted in Theorem 21.4, the result is valid also for half-integral $k$. The proof is the same as above; indeed, our discussions of §§21.1~8 include that case.
22.10. To show that the bound $s_{0}$ is best possible, we denote by $\mathcal{T}_{l}$ the vector space consisting of the functions $g$ of the form

$$
\begin{equation*}
g(z)=\sum_{\xi \in V} \lambda(\xi) \operatorname{det}(\xi)^{\mu \rho} \mathbf{e}_{\mathbf{a}}^{n}\left(\xi^{*} \tau \xi z\right) \quad\left(\lambda \in \mathcal{S}\left(V_{\mathbf{h}}\right), z \in \mathcal{H}\right) \tag{22.20}
\end{equation*}
$$

Here $V=K_{n}^{n}, \mu \in \mathbf{Z}^{\mathbf{b}} ; l=\mu+(n / 2) \mathbf{a}$ in Case SP and $l=\mu+n \mathbf{a}$ in Case UT; we assume $\mu_{v} \geq 0$ and $\mu_{v} \mu_{v \rho}=0$ for every $v \in \mathbf{a}$ in Case UT and $0 \leq \mu_{v} \leq 1$ for every $v \in \mathbf{a}$ in Case SP. As will be shown in $\S \S A 3.16$ and A5.5, $\mathcal{T}_{l}$ is contained in $\mathcal{M}_{l}$ and stable under $g \mapsto g \|_{l} \alpha$ for $\alpha \in G$. Let $\mathcal{T}_{l}^{\prime}=\mathcal{T}_{l} \cap \mathcal{S}_{l}$ and let $\mathcal{U}_{l}$ denote the orthogonal complement of $\mathcal{T}_{l}{ }^{\prime}$ in $\mathcal{S}_{l}$. Then we see that both $\mathcal{T}_{l}^{\prime}$ and $\mathcal{U}_{l}$ satisfy [S97, (10.7.1, 2, 3)]. Therefore we can speak of $\mathcal{T}_{l}^{\prime}(C)$ and $\mathcal{U}_{l}(C)$ for a subgroup $C$ of $G_{\mathbf{A}}$ of the above type. In fact, $\mathbf{g}$ of $\S 22.6$ belongs to $\mathcal{T}_{l}\left(C^{\prime}, \psi^{\prime}\right)$, as its construction in $\S A 5.5$ shows. As noted in [S97, §10.9], for every $s \in G_{\mathbf{h}}$ the map $\mathbf{f} \mapsto \mathbf{f}(x s)$ sends $\mathcal{T}_{l}^{\prime}(C)$ and $\mathcal{U}_{l}(C)$ onto $\mathcal{T}_{l}^{\prime}\left(s C s^{-1}\right)$ and $\mathcal{U}_{l}\left(s C s^{-1}\right)$. Now we have an inner product on $\mathcal{S}_{l}(C)$ defined by [S97, (10.9.6)]. Then $\mathcal{U}_{l}(C)$ is the orthogonal complement of $\mathcal{T}_{l}^{\prime}(C)$ in $\mathcal{S}_{l}(C)$. Indeed, that $\mathcal{T}_{l}^{\prime}(C)$ and $\mathcal{U}_{l}(C)$ are orthogonal is obvious. Suppose $\mathbf{f} \in \mathcal{S}_{l}(C)$ and $\left\langle\mathbf{f}, \mathcal{T}_{l}^{\prime}(C)\right\rangle=0$. Given $g \in \mathcal{T}_{l}^{\prime}$, we can define $\mathbf{g} \in \mathcal{T}_{l}\left(C_{1}\right)$ with a sufficiently small open subgroup $C_{1}$ of $C$ so that $g$ is one of the components of $\mathbf{g}$. Put $C=\bigsqcup_{y \in Y} y C_{1}$ with $Y \subset G_{\mathbf{h}}$ and $\mathbf{g}_{1}(x)=\sum_{y \in Y} \mathbf{g}(x y)$. Then $\mathbf{g}_{1} \in \mathcal{T}_{l}^{\prime}(C)$ and $0=\left\langle\mathbf{f}, \mathbf{g}_{1}\right\rangle=\#(Y)\langle\mathbf{f}, \mathbf{g}\rangle$ by $[S 97,(10.9 .8)]$. From this we easily see that every component of $\mathbf{f}$ is orthogonal to $\mathcal{T}_{l}^{\prime}$, that is, $\mathbf{f} \in \mathcal{U}_{l}(C)$ as expecte 1 .
22.11. Theorem. Let $s_{0}=(3 n / 2)+2-[K: F]$; let $\mathbf{f}$ be an eigenform belonging to $\mathcal{S}_{k}(C, \psi)$ with integral or half-integral $k$ as in $\S 20.11$ and let $\chi$ be a Hecke character of $K$. Then $\mathcal{Z}(s, \mathbf{f}, \chi)$ has a pole at $s=s_{0}$ only if the components of $\mathbf{f}$ in the sense of (20.3b) belong to $\mathcal{T}_{l}^{\prime}$ with $l=\mu+(n / 2) \mathbf{a}$ or $l=\mu+n \mathbf{a}$ with
$\mu \in \mathbf{Z}^{\mathbf{b}}$ as in §22.10. Moreover, for such $l$ and $\mu$ and an arbitrary $\chi$, there exists an eigenform $\mathbf{f} \in \mathcal{T}_{l}^{\prime}\left(C^{\prime}\right)$ for some $C^{\prime}$ such that $\mathcal{Z}(s, \mathbf{f}, \chi)$ has a pole at $s=s_{0}$, provided $\mu \neq 0$.

Proof. Since every Euler factor of $\mathcal{Z}(s, \mathbf{f}, \chi)$ is nonvanishing at $s_{0}$, we may remove any finite number of Euler factors by assuming $\mathfrak{c}=\mathfrak{e}$ and changing $\mathfrak{c}$ for its suitable multiple. Then $\mathbf{f} \in \mathcal{S}_{k}(C)$ and we may assume $\psi$ to be trivial; thus we have (22.19) with $\mathcal{Z}^{\prime}(s)=\mathcal{Z}(s, \mathbf{f}, \chi)$. We can even assume that $P=1$ with a suitable choice of $\mathfrak{c}$, since only finitely many ideals are involved in $P$. Then $\mathcal{Z}(s, \mathbf{f}, \chi)$ has a pole at $s=s_{0}$ if and only if $D(s, \mathbf{f}, \chi)$ has a pole at $s=s_{0}$, since $\Lambda$ is finite and nonzero at $s_{0}$. Now, by Proposition 22.2 (1), the right-hand side of (22.18b) has a pole at $s=0$ only if $k=l$ in Case SP and $k_{v}+k_{v \rho}=l_{v}+l_{v \rho}$ for every $v \in \mathbf{a}$ in Case UT. Assuming these conditions on $k$ and $l$, take $C$ so that $\operatorname{det}\left(G \cap p C p^{-1}\right)=1$ for every $p \in G_{\mathrm{h}}$; such a $C$ always exists in view of (4.33) and (4.34). Then we have $\mathcal{S}_{k}(C)=\mathcal{S}_{l}(C)$. In view of our remark about the map $\mathbf{f}(x) \mapsto \mathbf{f}(x s)$ in $\S 22.10$, both $\mathcal{T}_{l}^{\prime}(C)$ and $\mathcal{U}_{l}(C)$ are stable under $T(\mathfrak{a})$, and so, by Lemma $20.12(2), \mathcal{S}_{l}(C)$ is spanned by some eigenforms, each of which belongs to either $\mathcal{T}_{l}^{\prime}(C)$ or $\mathcal{U}_{l}(C)$. Suppose $\mathbf{f}$ is an eigenform belonging to $\mathcal{U}_{l}(C)$. Given $\chi$, define $\mathbf{g}$ as in $\S 22.6$. Then $\left\langle f_{p}, g_{p}\right\rangle=0$ for every $p$, and so by Proposition 22.2 (3), (22.18b), and (22.19), $\mathcal{Z}(s, \mathbf{f}, \chi)$ is finite at $s=s_{0}$. Thus $\mathbf{f} \in \mathcal{T}_{l}^{\prime}(C)$ if it has a pole at $s_{0}$.

Next, given $k$ as above with $\mu \neq 0$, take a Hecke character $\chi$ and define $\mathbf{g}$ as in $\S 22.6$ with $r=1_{n}$ and an arbitray $\tau \in S^{+}$. Take also an integral ideal $\mathfrak{a}$ in $F$ so that if $\xi \in G L_{n}(K) ; \xi^{*} \tau \xi=\tau$, and $\xi-1 \prec \mathfrak{r a}$, then $\xi=1$. Change $\mathfrak{c}$ for $\mathfrak{a c}$; then from (22.15) we see that $c_{\mathbf{g}}(\tau, 1) \neq 0$, so that $\mathbf{g} \neq 0$. We have $\mathbf{g} \in \mathcal{T}_{l}^{\prime}\left(C^{\prime}\right)$ with some $C^{\prime}$. Suppose $\mathcal{Z}(s, \mathbf{f}, \chi)$ is finite at $s=s_{0}$ for every eigenform $\mathbf{f} \in \mathcal{T}_{l}^{\prime}\left(C^{\prime}\right)$; then we see that $D\left(s, \mathbf{f}^{\prime}, \chi\right)$ is finite at $s=s_{0}$ for every $\mathbf{f}^{\prime} \in \mathcal{T}_{l}^{\prime}\left(C^{\prime}\right)$, which is not necessarily an eigenform. In particular $D(s, \mathbf{g}, \chi)$ is finite at $s=s_{0}$. This is a contradiction, since Proposition 22.2 (3) together with (22.18b) with $\mathbf{f}=\mathbf{g}$ shows that $D(s, \mathbf{g}, \chi)$ has a positive residue at $s_{0}$. This completes the proof.
22.12. Remark. (I) In the above we expressed $\mathcal{Z}(s, \mathbf{f}, \chi)$ times some factors as a finite linear combination of the functions $D_{q, \kappa}(s ; \mathbf{f}, \mathbf{g})$, each of which can be given by an integral of the form (22.9). In fact, at least if $\mathfrak{e}=\mathfrak{g}$ in Case SP, we can express $\mathcal{Z}(s, \mathbf{f}, \chi)$ itself as a single integral similar to that of (22.9) times some elementary factors. We refer the reader to [S96, (4.1)] and [S94a, (8.11)] for the explicit forms.
(II) If $G=S L_{2}(F)$, that is, if $n=1$ in Case SP, then we have a result supplementary to the above theorem. Indeed:
22.13. Theorem. We have $\mathcal{Z}(s, \mathbf{f}, \chi) \neq 0$ for $\operatorname{Re}(s) \geq 2$ in the following three cases: (i) $k=[k]$; (ii) $k \neq[k]$ and $k_{v}>3 / 2$ for some $v \in \mathbf{a}$; (iii) $k=3 \mathbf{a} / 2$ and $\mathbf{f}$ is orthogonal to $\mathcal{T}_{3 \mathbf{a} / 2}^{\prime}$. Moreover, in these cases $\mathcal{Z}(s, \mathbf{f}, \chi)$ is finite for every $s \in \mathbf{C}$ except that it has a possible simple pole at $s=1$ and $s=2$ if $k=[k]$.

If $k=[k]$, the function $\mathcal{Z}(s, \mathbf{f}, \chi)$ coincides with the function of [S91, (10.3)] as explained at the end of [S95a, §6], and hence nonvanishing follows from [S91, Proposition 3.3 and (10.4)]. The information concerning possible poles is given in [S95a, Theorem 6.4]. As for the case $k \neq[k]$, we refer the reader to [S95b, Proposition 6.2] and the succeeding paragraph. In this case $\mathcal{Z}$ is entire.

## CHAPTER VI

## ANALYTIC CONTINUATION AND NEAR HOLOMORPHY OF EISENSTEIN SERIES OF GENERAL TYPES

## 23. Eisenstein series of general types

23.1. To emphasize the dimensionality, denote the group $U\left(\eta_{n}\right)$ and the space $\mathcal{H}$ in Cases SP and UT by $G^{n}$ and $\mathcal{H}^{n}$, and the symbols $\mathcal{M}_{k}$ and $\mathcal{S}_{k}$ by $\mathcal{M}_{k}^{n}$ and $\mathcal{S}_{k}^{n}$; we write $M_{\mathbf{A}}^{n}$ and $\mathfrak{M}^{n}$ for $M_{\mathbf{A}}$ and $\mathfrak{M}$. We make the convention that $G^{0}=\{1\}$ and $\mathcal{M}_{k}^{0}=\mathcal{S}_{k}^{0}=\mathbf{C}$; also we understand that pr is the identity map of $G_{\mathbf{A}}^{n}$ onto itself if the weight is integral. We are going to introduce Eisenstein series associated with elements of $\mathcal{S}_{k}^{r}$ on subgroups $G^{r}$ embedded in $G^{n}$. These include those of Section 16 as special cases. The series of this type in the unitary case was treated in [S97] and we indicated there that the symplectic case can be handled in the same manner. In this section we present a more detailed treatment in Cases SP and UT, which, in Case UT, is essentially included in what was done in [S97]. Fixing an integer $r$ such that $0 \leq r \leq n$, we write each element $\alpha$ of $\left(K_{\mathbf{A}}\right)_{2 n}^{2 n}$ in the form

$$
\alpha=\left[\begin{array}{llll}
a_{1} & a_{2} & b_{1} & b_{2}  \tag{23.1}\\
a_{3} & a_{4} & b_{3} & b_{4} \\
c_{1} & c_{2} & d_{1} & d_{2} \\
c_{3} & c_{4} & d_{3} & d_{4}
\end{array}\right],
$$

where $x_{1}$ is of size $r$, and $x_{4}$ is of size $n-r$. Then we write $x_{i}=x_{i}(\alpha)$ and $x_{\alpha}=x(\alpha)=\left[\begin{array}{ll}x_{1}(\alpha) & x_{2}(\alpha) \\ x_{3}(\alpha) & x_{4}(\alpha)\end{array}\right]$ for $x=a, b, c, d$, and $i=1,2,3,4$. We understand that $x(\alpha)=x_{1}(\alpha)$ if $r=n$, and $x(\alpha)=x_{4}(\alpha)$ if $r=0$. We define a parabolic subgroup $P^{n, r}$ of $G^{n}$ by

$$
\begin{array}{r}
P^{n, r}=\left\{\alpha \in G^{n} \mid a_{2}(\alpha)=c_{2}(\alpha)=0, c_{3}(\alpha)=d_{3}(\alpha)=0, c_{4}(\alpha)=0\right\} \\
(0<r<n), \tag{23.2b}
\end{array}
$$

and define also maps $\pi_{r}:\left(K_{\mathbf{A}}\right)_{2 n}^{2 n} \rightarrow\left(K_{\mathbf{A}}\right)_{2 r}^{2 r}$ and $\lambda_{r}:\left(K_{\mathbf{A}}\right)_{2 n}^{2 n} \rightarrow K_{\mathbf{A}}$ by

$$
\pi_{r}(\alpha)=\left[\begin{array}{ll}
a_{1}(\alpha) & b_{1}(\alpha)  \tag{23.3}\\
c_{1}(\alpha) & d_{1}(\alpha)
\end{array}\right], \quad \lambda_{r}(\alpha)=\operatorname{det}\left(d_{4}(\alpha)\right)
$$

These define homomorphisms $P_{\mathbf{A}}^{n, r} \rightarrow G_{\mathbf{A}}^{r}$ and $P_{\mathbf{A}}^{n, r} \rightarrow K_{\mathbf{A}}^{\times}$. Clearly $\lambda_{0}(\alpha)=$ $\operatorname{det}\left(d_{\alpha}\right)$. We understand that $G_{\mathbf{A}}^{0}=\pi_{0}\left(P_{\mathbf{A}}^{n, 0}\right)=1, \pi_{n}(\alpha)=\alpha$, and $\lambda_{n}(\alpha)=1$ for $\alpha \in G_{\mathbf{A}}^{n}$. Notice that $P^{n, r}$ is a parabolic subgroup of $G^{n}$ in the sense of [S97, Section 2]; see [S97, §2.11] in particular.

Assuming $r>0$, for $z \in \mathbf{C}_{n}^{n}$ we let $\wp_{r}(z)$ denote the upper left $(r \times r)$-block of $z$, and use the same letter $\wp_{r}$ for the map $\left(\mathbf{C}_{n}^{n}\right)^{\mathbf{a}} \rightarrow\left(\mathbf{C}_{r}^{r}\right)^{\mathbf{a}}$ defined by $\wp_{r}\left(\left(z_{v}\right)_{v \in \mathbf{a}}\right)=$ $\left(\wp_{r}\left(z_{v}\right)_{v \in \mathbf{a}}\right)$. Then for $\alpha \in P_{\mathbf{A}}^{n, r}$ and $z \in \mathcal{H}^{n}$ we have

$$
\begin{equation*}
\wp_{r}(\alpha z)=\pi_{r}(\alpha) \wp_{r}(z), \quad j(\alpha, z)=\lambda_{r}\left(\alpha_{\mathbf{a}}\right) j\left(\pi_{r}(\alpha), \wp_{r}(z)\right) \tag{23.4}
\end{equation*}
$$

Here $\lambda_{r}\left(\alpha_{\mathbf{a}}\right)=\left(\lambda_{r}(\alpha)_{v}^{\rho}, \lambda_{r}(\alpha)_{v}\right)_{v \in \mathbf{a}} \in \mathbf{C}^{\mathbf{b}}$ in Case UT. For $\beta \in G_{\mathbf{A}}^{r}$ and $\gamma \in G_{\mathbf{A}}^{n-r}$ we define an element $\beta \times \gamma$ of $G_{\mathbf{A}}^{n}$ by

$$
\beta \times \gamma=\left[\begin{array}{cccc}
a_{\beta} & 0 & b_{\beta} & 0  \tag{23.5}\\
0 & a_{\gamma} & 0 & b_{\gamma} \\
c_{\beta} & 0 & d_{\beta} & 0 \\
0 & c_{\gamma} & 0 & d_{\gamma}
\end{array}\right] .
$$

Writing $\mathcal{G}^{n}$ for the group $\mathcal{G}$ defined as in $\S 14.14$, we put

$$
\begin{equation*}
\mathcal{P}^{n, r}=\left\{(\alpha, p) \in \mathcal{G}^{n} \mid \alpha \in P^{n, r}\right\} \tag{23.6}
\end{equation*}
$$

and define homomorphisms $\pi_{r}: \mathcal{P}^{n, r} \rightarrow \mathcal{G}^{r}$ and $\lambda_{r}: \mathcal{P}^{n, r} \rightarrow F_{\mathbf{a}}^{\times}$by

$$
\begin{gather*}
\pi_{r}((\alpha, p))=\left(\pi_{r}(\alpha),\left|\lambda_{r}(\alpha)\right|^{-\mathbf{a} / 2} p^{\prime}\right) \quad \text { with } \quad p^{\prime}(z)=p\left(\left[\begin{array}{cc}
z & w \\
t^{t} w & z^{\prime}
\end{array}\right]\right),  \tag{23.7a}\\
\lambda_{r}((\alpha, p))=\lambda_{r}\left(\alpha_{\mathbf{a}}\right) \tag{23.7b}
\end{gather*}
$$

where it should be observed that $p^{\prime}(z)$ does not depend on the choice of $w$ and $z^{\prime}$.
23.2. We now fix a weight $k$ as in Section 16, and make the convention that $\mathcal{G}^{n}=G^{n}$ and $\mathcal{P}^{n, r}=P^{n, r}$ if $k$ is integral. We put $[k]=k$ if $k$ is integral, and $[k]=\left(k_{v}-1 / 2\right)_{v \in \mathbf{a}}$ otherwise; we also put $m=k$ and $\ell=[k]$ in Case SP, and $m=\left(k_{v \rho}+k_{v}\right)_{v \in \mathbf{a}}$ and $\ell=\left(k_{v}-k_{v \rho}\right)_{v \in \mathbf{a}}$ in Case UT. Clearly

$$
\begin{equation*}
j^{k}(\xi, z)=j^{k}\left(\pi_{r}(\xi), \wp_{r}(z)\right) \lambda_{r}(\xi)^{[k]}\left|\lambda_{r}(\xi)\right|^{k-[k]} \quad \text { if } \quad \xi \in \mathcal{P}^{n, r} \tag{23.8}
\end{equation*}
$$

For $0 \leq r \leq n$ and a congruence subgroup $\Gamma$ of $\mathcal{G}^{n}$ we put

$$
\begin{align*}
& \mathcal{M}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}\right)  \tag{23.9}\\
& =\left\{\left.f \in \mathcal{M}_{k}^{r}\left|f \|_{k} \pi_{r}(\gamma)=\lambda_{r}(\gamma)^{\ell}\right| \lambda_{r}(\gamma)\right|^{-\ell} f \quad \text { for every } \quad \gamma \in \Gamma \cap \mathcal{P}^{n, r}\right\} \\
& \mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}\right)=\mathcal{M}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}\right) \cap \mathcal{S}_{k}^{r} \tag{23.9a}
\end{align*}
$$

These spaces are $\{0\}$ unless the following condition is satisfied:
(23.10) There is a homomorphism $\varphi$ of $\pi_{r}\left(\Gamma \cap \mathcal{P}^{n, r}\right)$ into $\mathbf{T}$ such that $\varphi\left(\pi_{r}(\gamma)\right)=$ $\lambda_{r}(\gamma)^{\ell}\left|\lambda_{r}(\gamma)\right|^{-\ell}$ for every $\gamma \in \Gamma \cap \mathcal{P}^{n, r}$.
Clearly $\varphi^{2}=1$ in Case SP. As will be shown in the proof of Lemma 23.13 below, $\varphi$ is of finite order also in Case UT. Under (23.10), $\mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}\right)$ consists of all $f \in \mathcal{S}_{k}^{r}$ such that $f \|_{k} \varepsilon=\varphi(\varepsilon) f$ for every $\varepsilon \in \pi_{r}\left(\Gamma \cap \mathcal{P}^{n, r}\right)$. For $f \in \mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}\right), z \in \mathcal{H}_{n}$, and $s \in \mathbf{C}$ we put

$$
\begin{equation*}
\delta(z, s ; f)=f\left(\wp_{r}(z)\right)\left[\delta(z) / \delta\left(\wp_{r}(z)\right)\right]^{s \mathbf{a}-m / 2} \tag{23.11}
\end{equation*}
$$

where $\delta(z)=\left(\operatorname{det}\left((i / 2)\left(z^{*}-z\right)\right)_{v}\right)_{v \in \mathbf{a}}$; we understand that $\delta\left(\wp_{0}(z)\right)=1$. We note that

$$
\begin{equation*}
\delta(z, s ; f) \|_{k} \beta=\delta\left(z, s ; f \|_{k} \pi_{r}(\beta)\right) \lambda_{r}(\beta)^{-\ell}\left|\lambda_{r}(\beta)\right|^{\ell-2 s \mathbf{a}} \quad \text { if } \quad \beta \in \mathcal{P}^{n, r} \tag{23.12}
\end{equation*}
$$

and slso that $a^{-\ell}|a|^{\ell}=a^{-k}|a|^{m}$ for every $a \in K^{\times}$if $k \in \mathbf{Z}^{\mathbf{b}}$.
We now define an Eisenstein series $E_{k}^{n, r}(z, s ; f, \Gamma)$ by

$$
\begin{equation*}
E_{k}^{n, r}(z, s ; f, \Gamma)=\sum_{\gamma \in A} \delta(z, s ; f) \|_{k} \gamma, \quad A=\left(\Gamma \cap \mathcal{P}^{n, r}\right) \backslash \Gamma \tag{23.13}
\end{equation*}
$$

The sum is formally well-defined, since $f \in \mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}\right)$. It is convergent for $\operatorname{Re}(2 s)>n+r+1$ in Case SP and $\operatorname{Re}(s)>n+r$ in Case UT (see [S97, Proposition A3.7 and $\S A 3.9])$. For $\Gamma^{\prime} \subset \Gamma$ we can easily verify that

$$
\begin{equation*}
\left[\Gamma \cap P: \Gamma^{\prime} \cap P\right] E_{k}^{n, r}(z, s ; f, \Gamma)=\sum_{\alpha \in \Gamma^{\prime} \backslash \Gamma} E_{k}^{n, r}\left(z, s ; f, \Gamma^{\prime}\right) \|_{k} \alpha \tag{23.13a}
\end{equation*}
$$

The series $E_{k}^{n, r}(z, s ; f, \Gamma)$ can be defined even when $r=0$ or $n$. If $r=n$, we have $\mathcal{S}_{k}^{n}\left(\Gamma, \mathcal{P}^{n, n}\right)=\mathcal{S}_{k}^{n}(\Gamma)$ and

$$
\begin{equation*}
E_{k}^{n, n}(z, s ; f, \Gamma)=f \tag{23.14}
\end{equation*}
$$

If $r=0$, we have $\delta(z, s ; c)=c \delta(z)^{s \mathbf{a}-m / 2}$ for a constant $c \in \mathbf{C}$, and

$$
\mathcal{S}_{k}^{0}\left(\Gamma, \mathcal{P}^{n, 0}\right)=\left\{\begin{array}{lcc}
\mathbf{C} & \text { if } & \lambda_{0}(\gamma)^{-\ell}\left|\lambda_{0}(\gamma)\right|^{\ell}=p_{\gamma} \quad \text { for every } \quad \gamma \in \Gamma \cap \mathcal{P}^{n, 0}  \tag{23.15}\\
\{0\} & \text { otherwise }
\end{array}\right.
$$

where we understand that $\gamma=\left(\operatorname{pr}(\gamma), p_{\gamma}\right)$ if $k \neq[k]$ and $p_{\gamma}=1$ if $k=[k]$. Thus

$$
\begin{equation*}
E_{k}^{n, 0}(z, s ; c, \Gamma)=c \sum_{\gamma \in A} \delta^{s \mathbf{a}-m / 2} \|_{k} \gamma, \quad A=\left(\Gamma \cap \mathcal{P}^{n, 0}\right) \backslash \Gamma \tag{23.16}
\end{equation*}
$$

We can now ask questions similar to ( $\mathrm{R} 1,2,3$ ) of $\S 17.3$ about the nature of $E_{k}^{n, r}(z, s ; f, \Gamma)$ for some values of $s$. The answer will given in Theorem 23.11.
23.3. To define the adelized version of our Eisenstein series, we take a set of data ( $k, \mathfrak{b}, \mathfrak{c}, \mathfrak{e}, \chi$ ) as in $\S \S 16.5$ and 19.1 ; we assume conditions ( $16.24 \mathrm{a}, \mathrm{b}, \mathrm{c}$ ). We denote the subgroup $D[\mathfrak{x}, \mathfrak{y}]$ of $G_{\mathbf{A}}^{n}$ by $D^{n}\{\mathfrak{x}, \mathfrak{y}]$ and similarly $C$ of (19.1) by $C^{n}$. For simplicity let us put $D^{n}=D^{n}\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right], D_{0}^{n}=D_{0}^{n}\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]$, and

$$
\begin{align*}
C^{n, r} & =\left\{x \in D^{n} \mid a_{1}(x)-1 \prec \mathfrak{r e}, a_{2}(x) \prec \mathfrak{r e}, b_{1}(x) \prec \mathfrak{r b}^{-1} \mathfrak{e}\right\},  \tag{23.17a}\\
C_{0}^{n, r} & =D_{0}^{n} \cap C^{n, r} . \tag{23.17b}
\end{align*}
$$

It can easily be verified that $C^{n, r}$ is indeed a subgroup of $G_{\mathbf{A}}^{n}$ and that $\pi_{r}\left(P^{n, r} \cap\right.$ $\left.C^{n, r}\right) \subset C^{r}$. If $x=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right] \in D^{n}$, we see that $a d^{*}-1 \prec \mathfrak{r c}$, and hence the first two conditions on $a_{1}$ and $a_{2}$ can be replaced by $d_{1}(x)-1 \prec \mathfrak{r e}$ and $d_{3}(x) \prec \mathfrak{r e}$.

For every $g \in G_{\mathbf{h}}^{n}$ we have

$$
\begin{equation*}
G_{\mathbf{A}}^{n}=P_{\mathbf{A}}^{n, r} g D_{0}^{n}\left[\mathfrak{b}^{-1}, \mathfrak{b}\right] g^{-1} \tag{23.18}
\end{equation*}
$$

We already noted this in (16.22) when $g=1$ and $r=0$, but the equality with $g=1$ is true for an arbitrary $r$ by virtue of [S97, Propositions 6.13, 7.2, and 7.12]. Then the case of more general $g$ follows from that special case by an easy principle; see a remark at the end of [S97, Section 5]. Now, taking $g \in P_{\mathrm{h}}^{n, r}$, we define an $\mathfrak{r}$-ideal $\mathfrak{a}_{r}^{g}(x)$ for $x \in G_{\mathbf{A}}^{n}$ and also an $\mathbf{R}$-valued function $\varepsilon_{r}$ on $G_{\mathbf{A}}^{n}$ by

$$
\begin{align*}
& \mathfrak{a}_{r}^{g}(x)=\lambda_{r}(p) \mathfrak{r} \quad \text { if } x \in p g D^{n}\left[\mathfrak{b}^{-1}, \mathfrak{b}\right] g^{-1} \text { with } p \in P_{\mathbf{A}}^{n, r}  \tag{23.19}\\
& \varepsilon_{r}(x)=\left|\lambda_{r}(p) \lambda_{r}(p)^{\rho}\right|_{F} \quad \text { if } \quad x \in p D_{0}^{n}\left[\mathfrak{b}^{-1}, \mathfrak{b}\right] \text { with } p \in P_{\mathbf{A}}^{n, r} .
\end{align*}
$$

These are well-defined, and depend on $\mathfrak{b}$. If $r=0$ and $g=1, \mathfrak{a}_{r}^{g}(x)$ coincides with $\mathrm{il}_{\mathfrak{b}}(x)$ of (16.23); also $\varepsilon_{0}$ coincides with $\varepsilon$ of (16.23). We have

$$
\begin{align*}
& \varepsilon_{r}\left(x_{\mathbf{h}}\right)=N\left(\mathfrak{a}_{r}^{1}(x)\right)^{-u}, \text { where } u=2 /[K: F],  \tag{23.21a}\\
& \varepsilon_{r}\left(x_{\mathbf{a}}\right)=\left[\delta(z) / \delta\left(\wp_{r}(z)\right)\right]^{-\mathbf{a}} \text { if } z=x(\mathbf{i}) . \tag{23.21b}
\end{align*}
$$

The first equality is obvious; the latter follows easily from (23.4). For $x \in M_{\mathbf{A}}^{n}$ we put $\mathfrak{a}_{r}^{g}(x)=\mathfrak{a}_{r}^{g}(\operatorname{pr}(x))$ and $\varepsilon_{r}(x)=\varepsilon_{r}(\operatorname{pr}(x))$.

In Case UT our treatment in this section is a special case of what was done in [S97, §§12.5~12.9]. This can be seen by taking ( $G^{n}, D\left[\mathfrak{b}, \mathfrak{b}^{-1}\right], C^{n, r} ; G^{r}, C^{r}$; $P^{n, r}, \pi_{r}, \lambda_{r}$ ) here to be ( $G^{\psi}, C^{\psi}, D^{\psi} ; G^{\varphi}, D^{\varphi} ; P, \pi, \lambda_{0}$ ) there; define also $\xi^{\psi}$ and $\xi^{\varphi}$ there by $\xi^{\psi}(w)=\chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{w}\right)\right)^{-1}$ for $w \in C^{n, r}$ and $\xi^{\varphi}(x)=\chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{x}\right)\right)^{-1}$ for $x \in C^{r}$. Case SP can be handled in the same manner as noted in [S97, §12.13]; a similar but a somewhat different treatment was given in [S95a].

Now we consider $\mathcal{S}_{k}^{r}\left(C^{r}, \chi_{0}\right)$ as defined in $\S 20.1$ with $\chi_{0}$ given by $\chi_{0}(a)=$ $\chi\left(a^{\rho}\right)^{-1}$. Taking an element $\mathbf{f}$ of $\mathcal{S}_{k}^{r}\left(C^{r}, \chi_{0}\right)$, we define a function $\mu$ on $G_{\mathbf{A}}^{n}$ or on $M_{\mathbf{A}}^{n}$ according as $k$ is integral or half-integral as follows: $\mu(x)=0$ if $\operatorname{pr}(x) \notin$ $P_{\mathbf{A}}^{n, r} C^{n, r}$; if $x=p w$ with $\operatorname{pr}(p) \in P_{\mathbf{A}}^{n, r}$ and $\operatorname{pr}(w) \in C_{0}^{n, r}$, then

$$
\begin{equation*}
\mu(x)=\chi\left(\lambda_{r}(p)\right)^{-1} \chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{w}\right)\right)^{-1} j_{w}^{k}(\mathbf{i})^{-1} \mathbf{f}\left(\pi_{r}(p)\right) \tag{23.22}
\end{equation*}
$$

This is well-defined. Notice that $b^{k}=b^{\ell}$ if $b \in K_{\mathrm{a}}$ and $\left|b_{v}\right|=1$ for every $v \in \mathbf{a}$ in Case UT. For half-integral $k$ we have to establish $\pi_{r}(p)$ as an element of $M_{\mathbf{A}}^{r}$, which will be done in $\S 23.8$. We can easily verify that

$$
\begin{equation*}
\mu(\alpha x w)=\chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{w}\right)\right)^{-1} j_{w}^{k}(\mathbf{i})^{-1} \mu(x) \quad \text { if } \alpha \in P^{n, r} \text { and } \operatorname{pr}(w) \in C_{0}^{n, r} \tag{23.23}
\end{equation*}
$$

Then we define a function $E_{\mathbf{A}}\left(x, s ; \mathbf{f}, \chi, C^{n, r}\right)$ for $s \in \mathbf{C}$ and $x \in G_{\mathbf{A}}^{n}$ or $x \in M_{\mathbf{A}}^{n}$ according as $k$ is integral or half-integral by

$$
\begin{equation*}
E_{\mathbf{A}}(x, s)=E_{\mathbf{A}}\left(x, s ; \mathbf{f}, \chi, C^{n, r}\right)=\sum_{\alpha \in A} \mu(\alpha x) \varepsilon_{r}(\alpha x)^{-s}, \quad A=P^{n, r} \backslash G^{n} \tag{23.24}
\end{equation*}
$$

This is well-defined. The domain of convergence for (23.24) is the same as that for (23.13), as explained in [S97, §12.11].
23.4. Fix an element $g \in G_{\mathbf{h}}^{r}$ such that $g_{v}=1$ for every $v \mid$ c. By [S97, Lemma $9.8(3)]$ or by strong approximation on $S p(n, F)$, we can find a finite subset $\{q\}$ of $P_{\mathrm{h}}^{n, r}$ such that $G_{\mathbf{A}}^{n}=\bigsqcup_{q} G^{n} q C^{n, r}$ and

$$
\begin{equation*}
q=g \times \operatorname{diag}[\widehat{\varphi}, \varphi] \text { with } \varphi \in G L_{n-r}(K)_{\mathbf{h}}, q_{v}=1 \text { if } v \mid \mathbf{c} . \tag{23.25}
\end{equation*}
$$

We can take $g=1$, but for some technical reasons, we consider $q$ with $g$ of a more general type. For such a $q$ we have

$$
\begin{equation*}
P_{\mathbf{A}}^{n, r} q C^{n, r} q^{-1}=P_{\mathbf{A}}^{n, r} C^{n, r} \tag{23.25a}
\end{equation*}
$$

Indeed, if $v \nmid c$, then $P_{v}^{n, r} C_{v}^{n, r}=P_{v}^{n, r} D_{v}^{n}=G_{v}^{n}$, and hence $P_{v}^{n, r} q_{v} C_{v}^{n, r} q_{v}^{-1}=G_{v}^{n}$. Since $q_{v}=1$ for $v \mid \mathfrak{c}$, we obtain (23.25a).

Since $E_{\mathbf{A}}(x, s)$ satisfies a formula of type (20.3a), the principle of [S97, (10.7.5)] defines a function $E_{q}(z)=E_{q}(z, s ; \mathbf{f}, \chi, C)$ for $z \in \mathcal{H}^{n}$ for each $q$ by

$$
\begin{equation*}
E_{\mathbf{A}}(q y, s)=\left(E_{q} \|_{k} y\right)(\mathbf{i}) \quad \text { if } \operatorname{pr}(y) \in G_{\mathbf{a}}^{n} . \tag{23.26}
\end{equation*}
$$

The function $E_{\mathbf{A}}$ is completely determined by the functions $E_{q}$. In Case SP, $E_{1}$ determines $E_{\mathbf{A}}$.

Next, for each fixed $q$ we consider a complete set of representatives $\mathcal{R}_{q}$ for $P^{n, r} \backslash\left[G^{n} \cap P_{\mathrm{A}}^{n, r} q C^{n, r} q^{-1}\right]$. By [S97, Lemma 9.6(3,4) and Lemma 9.8(1)] we can find a finite subset $\{\zeta\}$ of $P_{\mathbf{h}}^{n, r}$ such that $G^{n} \cap \zeta q C^{n, r} q^{-1} \neq \emptyset, G^{n} \cap P_{\mathbf{A}}^{n, r} q C^{n, r} q^{-1}=$ $\bigsqcup_{\zeta}\left(G^{n} \cap P^{n, r} \zeta q C^{n, r} q^{-1}\right)$, and each $\zeta$ is of the form

$$
\begin{equation*}
\zeta=e \times \operatorname{diag}[\widehat{s}, s] \text { with } s \in G L_{n-r}(K)_{\mathbf{h}} \text { and } e \in \mathcal{B}^{r}, \zeta_{v}=1 \text { if } v \mid \mathfrak{c}, \tag{23.27}
\end{equation*}
$$

where $\mathcal{B}^{r}$ is any fixed subset of $G_{\mathrm{h}}^{r}$ satisfying (20.5) with $r$ in place of $n$. We can and do take $\mathcal{B}^{r}=\{1\}$ in Case SP . Then we can take $\mathcal{R}_{q}$ to be a subset of $\bigsqcup_{\zeta}\left(G^{n} \cap \zeta q C^{n . r} q^{-1}\right)$. With $\kappa$ as in (16.24a), we put

$$
\begin{equation*}
\delta(z, s ; f, \kappa)=f\left(\wp_{r}(z)\right)\left[\delta(z) / \delta\left(\wp_{r}(z)\right)\right]^{\mathbf{a} \mathbf{a}-(m-i \kappa) / 2} \tag{23.28}
\end{equation*}
$$

If $\kappa=0$, this is the same as (23.11).
23.5. Proposition. Suppose $\mathfrak{e}=\mathfrak{c}$; let $q$ and $g$ be as in (23.25); suppose $q=1$ in Case SP. Then we have

$$
\begin{aligned}
E_{q}(z)= & \chi_{\mathbf{h}}\left(\lambda_{r}(q)^{-1}\right)\left|\lambda_{r}(q)\right|_{K}^{-s} \\
& \cdot \sum_{\alpha \in \mathcal{R}_{q}} N\left(\mathfrak{a}_{r}^{q}(\alpha)\right)^{u s} \chi_{\mathbf{a}}\left(\lambda_{r}(\alpha)\right) \chi^{*}\left(\lambda_{r}(\alpha) \mathfrak{a}_{r}^{q}(\alpha)^{-1}\right) \delta\left(z, s ; f_{e g}, \kappa\right) \|_{k} \alpha
\end{aligned}
$$

where $u$ is as in (23.21a), $f_{\text {eg }}$ is the eg-component of $\mathbf{f}$ in the sense of (20.3b), and $e$ is the element of $G_{\mathrm{h}}^{r}$ in (23.27) for $\zeta$ such that $\alpha \in \zeta q C^{n, r} q^{-1}$.

In Case UT this is included in [S97, Proposition 12.10] as a special case. The proof there can easily be translated into the symplectic case. However, there is one nontrivial technical point in the case of half-integral $k$, which will be explained in §23.8.
23.6. Proposition. Let $X$ be the set of all Hecke characters of $K$ satisfying (16.24a, b) with $\kappa=0$. Suppose $\mathfrak{e}=\mathfrak{c}$ and $X \neq \varnothing$; put $\Gamma^{\prime}=\left\{\alpha \in G^{n} \cap\right.$ $\left.D\left[\mathfrak{b}^{-1} \mathfrak{c}, \mathfrak{b c}\right] \mid a_{\alpha}-1 \prec \mathfrak{c}\right\}$. Let $f \in \mathcal{S}_{k}^{r}\left(\Gamma, P^{n, r}\right)$ with a congruence subgroup $\Gamma$ of $\mathcal{G}^{n}$ containing $\Gamma^{\prime}$. Take $\mathcal{B}^{r}$ as in §23.4 with 1 as one of its members, and take $\mathbf{f} \in \mathcal{S}_{k}^{r}\left(C^{r}\right)$ so that $\mathbf{f} \leftrightarrow\left(f_{b}\right)_{b \in \mathcal{B}^{r}}$ with $f_{1}=f$ and $f_{b}=0$ for $1 \neq b \in \mathcal{B}^{r}$. Then

$$
\left[P^{n, r} \cap \Gamma: P^{n, r} \cap \Gamma^{\prime}\right] \#(X) E_{k}^{n, r}(z, s ; f, \Gamma)=\sum_{\chi \in X} \sum_{\xi \in \Gamma^{\prime} \backslash \Gamma} E_{1}\left(z, s ; \mathbf{f}, \chi, C^{n, r}\right) \| \xi
$$

where $E_{1}(\cdots)$ denotes $E_{q}(z)$ of (23.26) with $q=1$.
Here if $k$ is half-integral, we identify $\Gamma^{\prime}$ with its image in $\mathcal{G}^{n}$ under the map $\gamma \mapsto\left(\gamma, h_{\gamma}\right)$, which is meaningful, since $\Gamma^{\prime} \subset \Gamma^{\theta}$ by (16.24c).

The proof in Case UT was given in [S97, Proposition 20.10]. The symplectic case can be proven in a similar and simpler way; in fact it is an easy modification of the proof of Lemma 17.2 (2). We insert here an easy fact:
23.7. Lemma. Let $P^{\prime}=\bigcap_{r=0}^{n} P^{n, r}$. Then $P_{\mathbf{A}}^{n} D^{n}=P_{\mathbf{A}}^{\prime} D^{n}$.

Proof. Clearly it is sufficient to prove that $P_{\mathbf{A}}^{n} \subset P_{\mathbf{A}}^{\prime} D^{n}$. Let $\alpha \in P_{\mathbf{A}}^{n}$. By [S97, Proposition 3.5] we can find an element $g$ of $\prod_{v \in \mathbf{h}} G L_{n}\left(\mathfrak{r}_{v}\right)$ such that $\left(d_{\alpha} g^{-1}\right)_{\mathbf{h}}$ is upper triangular. Put $\gamma=\operatorname{diag}[\hat{g}, g]$. Then $\gamma \in D^{n}$ and $\left(\alpha \gamma^{-1}\right)_{\mathbf{h}} \in P_{\mathbf{h}}^{\prime}$, from which we obtain the desired fact.
23.8. The symbol $\pi_{r}(p)$ in (23.22) is meaningful as an element of $M_{\mathbf{A}}^{r}$, since there exists a homomorphism $\pi_{r}: \mathrm{pr}^{-1}\left(P_{\mathbf{A}}^{n . r}\right) \rightarrow M_{\mathbf{A}}^{r}$ with the following properties:
(23.29a) $\mathrm{pr} \circ \pi_{r}=\pi_{r} \circ \mathrm{pr}$.
(23.29b) If $\xi \in \operatorname{pr}^{-1}\left(P_{\mathbf{A}}^{n . r}\right)$ and $\pi_{r}(\xi) \in \mathfrak{M}^{r}$, then $\xi \in \mathfrak{M}^{n}$ and $j^{k}(\xi, z)=\lambda_{r}(\xi)^{[k]}$ $\cdot\left|\lambda_{r}(\xi)\right|^{k-[k]} j^{k}\left(\pi_{r}(\xi), \wp_{r}(z)\right)$.
(23.29c) The restrictions of $\pi_{r}$ and $\lambda_{r}$ to $\mathcal{P}^{n, r}$ coincide with the maps of (23. $7 \mathrm{a}, \mathrm{b}$ ).
(23.29d) $\ell_{r} \circ \pi_{r}=\pi_{r} \circ \ell_{n}$, where $\ell_{r}$ denotes the canonical lift $G^{r} \rightarrow M_{\mathbf{A}}^{r}$.

For the proof, see [S95a, Lemmas 3.4 and 3.5].
Now as to the proof of Proposition 23.5 for half-integral $k$, we take $y \in \operatorname{pr}^{-1}\left(G_{\mathbf{a}}^{n}\right)$ and put $z=y(\mathbf{i})$. Then $E_{1}(z)=\sum_{\alpha \in A} \mu(\alpha y) \varepsilon(\alpha y)^{-s} j_{y}^{k}(\mathbf{i})$. Since $\mu(\alpha y) \neq 0$ only if $\alpha \in G^{n} \cap P_{\mathbf{A}}^{n, r} C^{n, r}$, we can take $\alpha$ in $\mathcal{R}_{1}$. Then $\alpha \in \zeta C^{n, r}$ with some $\zeta$ as in (23.27). Clearly $\zeta \in P_{\mathbf{h}}^{n, 0} \cap P_{\mathbf{h}}^{n, r}$. Now we can put $\alpha_{\mathbf{a}} \operatorname{pr}(y)=x x^{\prime}$ with $x \in P_{\mathbf{a}}^{n, 0} \cap P_{\mathbf{a}}^{n, r}$ and $x^{\prime} \in G_{\mathbf{a}}^{n}$ such that $x^{\prime}(\mathbf{i})=\mathbf{i}$. This is well-known, and proved in [S97, p.54, lines 8~9]. Then $\zeta x \in P_{\mathbf{A}}^{n, 0} \cap P_{\mathbf{A}}^{n, r}$. Put $\sigma=r_{P}(\zeta x)$ with $r_{P}$ of (16.15) and $w=\sigma^{-1} \alpha y$. Then $\operatorname{pr}(w) \in C_{0}^{n, r}$ and $\operatorname{pr}\left(\pi_{r}(\sigma)\right)_{\mathbf{h}}=\pi_{r}(\operatorname{pr}(\sigma))_{\mathbf{h}}=\pi_{r}(\zeta)=1$, and hence $\pi_{r}(\sigma) \in \mathfrak{M}^{r}$. By (16.16c) and (23.29b) we have $j_{\alpha}^{k}(z) j_{y}^{k}(\mathbf{i})=j_{\alpha y}^{k}(\mathbf{i})=j_{\sigma w}^{k}(\mathbf{i})=j_{\sigma}^{k}(\mathbf{i}) j_{w}^{k}(\mathbf{i})=$ $j^{k}\left(\pi_{r}(\sigma), \mathbf{i}\right) \lambda_{r}(\sigma)^{[k]}\left|\lambda_{r}(\sigma)\right|^{k-[k]} j_{w}^{k}(\mathbf{i})$. Once this is established, we can repeat what was done in [S97, p.98] to obtain the formula of Proposition 23.5
23.9. Theorem. The notation being the same as in §23.3, suppose that $\mathfrak{e}=\mathfrak{c}$ and $n>r>0$. Assuming that $\mathbf{f} \mid T(\mathfrak{a})=\lambda(\mathfrak{a}) \mathbf{f}$ for every $\mathfrak{a}$, define $\mathcal{Z}(s, \mathbf{f}, \chi)$ by (20.21) or (21.7). Put

$$
\begin{aligned}
& \mathcal{F}_{q}(z, s ; \mathbf{f}, \chi, C)=E_{q}(z, s ; \mathbf{f}, \chi, C) \\
& \cdot \mathcal{Z}(u s, \mathbf{f}, \chi) \cdot \begin{cases}\prod_{i=r+1}^{[(n+r) / 2]} L_{\mathfrak{c}}\left(4 s-2 i, \chi^{2}\right) & \left(\text { Case SP, } k \in \mathbf{Z}^{\mathbf{a}}\right), \\
\prod_{i=r+1}^{[(n+r+1) / 2]} L_{\mathfrak{c}}\left(4 s-2 i+1, \chi^{2}\right) & \left(\text { Case SP, } k \notin \mathbf{Z}^{\mathbf{a}}\right), \\
\prod_{i=2 r}^{n+r-1} L_{\mathfrak{c}}\left(2 s-i, \chi_{1} \theta^{i}\right) & \text { (Case UT) } \\
\Gamma_{k, \kappa}^{n, r}(s)=\prod_{v \in \mathbf{a}} \Gamma_{r}^{\iota}\left(s-\lambda_{r}+\left(m_{v}+i \kappa_{v}\right) / 2\right) \gamma\left(s+\left(i \kappa_{v} / 2\right), m_{v}\right)\end{cases}
\end{aligned}
$$

with $\gamma(s, a)$ defined in both cases as follows:
Case SP:

$$
\gamma(s, a)= \begin{cases}\Gamma\left(s+\frac{a}{2}-\left[\frac{2 a+n+r}{4}\right]\right) \Gamma_{n}^{1}\left(s+\frac{a-r}{2}\right) & \left(a \in \mathbf{Z}, a \geq \frac{n+r}{2} \in \mathbf{Z}\right), \\ \Gamma_{n}^{1}(s+(a-r) / 2) & (a \in \mathbf{Z}, a>(n+r) / 2 \notin \mathbf{Z}), \\ \Gamma_{2 a+1-r}^{1}\left(s+\frac{a-r}{2}\right) \prod_{i=a+1}^{[(n+r) / 2]} \Gamma(2 s-i) & \left(a \in \mathbf{Z}, \frac{r}{2} \leq a<\frac{n+r}{2}\right), \\ \Gamma_{n}^{1}(s+(a-r) / 2) & (a \notin \mathbf{Z}, a>(n+r) / 2 \in \mathbf{Z}), \\ \Gamma\left(s+\frac{a-1}{2}-\left[\frac{2 a+n+r-2}{4}\right]\right) \Gamma_{n}^{1}\left(s+\frac{a-r}{2}\right) & \left(a \notin \mathbf{Z}, a>\frac{n+r}{2} \notin \mathbf{Z}\right), \\ \Gamma_{2 a+1-r}^{1}\left(s+\frac{a-r}{2}\right) \prod_{i=[a]+1}^{[(n+r-1) / 2]} \Gamma\left(2 s-i-\frac{1}{2}\right) & \left(a \notin \mathbf{Z}, \frac{r}{2} \leq a \leq \frac{n+r}{2}\right)\end{cases}
$$

Case UT: $\quad \gamma(s, a)=q^{n+r}(s, a) \Gamma_{n}^{2}(s-r+(a / 2))$.

Here $\iota=[K: F], u=2 / \iota, m=k$ in Case SP, $m=\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$ in Case UT; $\lambda_{r}=(r+1) / 2$ in Case SP and $\lambda_{r}=r$ in Case UT; $\theta$ and $\chi_{1}$ are the same as in (20.20); $q^{n+r}(s, a)$ is defined as in Theorem 20.14 (with $n+r$ in place of $t$ ). Then the product $\Gamma_{k, \kappa}^{n, r}(s) \mathcal{F}_{q}(z, s ; \mathbf{f}, \chi, C)$ can be continued to the whole s-plane as a meromorphic function with finitely many poles, which are all simple. The set of poles of this product is contained in the set of poles of the function $\mathcal{P}(s)$ defined as follows: in Case SP, $\mathcal{P}$ is the product of Theorem 16.11 defined with $\{n+r, k, \chi\}$ as $\{n, k, \chi\}$ there; in Case UT, $\mathcal{P}$ is the product given in [S97, Theorem 19.3] defined with $\{n+r, m, \chi\}$ as $\{n, k, \chi\}$ there.

The proof will be completed in $\S 25.6$. Notice that $k_{v} \geq n / 2$ in Case SP and $m_{v} \geq n$ in Case UT for every $v \in \mathbf{a}$ by (6.42).

In the above theorem we assumed $r>0$. If $r=0$ and $\mathbf{f}$ is constant 1 , then we easily see that the function of (23.24) is exactly that of (16.27), and the conclusion of Theorem 23.9 is reduced to Theorem 16.11 in Case SP and to [S97, Theorem 19.3] in Case UT, if we simply take $\mathcal{Z}(u s, f, \chi)$ to be 1 and $r=0$ in the definition of various objects.
23.10. Corollary. The function $E_{k}^{n, r}(z, s ; f, \Gamma)$ of (23.13) can be continued as a meromorphic function of $s$ to the whole complex plane. In particular, it is holomorphic for $\operatorname{Re}(4 s)>\operatorname{Max}(n+r+2,3 r+2)$ if $\operatorname{Max}_{v \in \mathbf{a}} 2\left|k_{v}\right| \geq n+r$ in Case $S P$, and for $\operatorname{Re}(2 s)>\operatorname{Max}(n+r, 3 r)$ if $\operatorname{Min}_{v \in \mathbf{a}}\left|k_{v}+k_{v \rho}\right| \geq n+r$ in Case UT.

Proof. Let $X$ and $\Gamma^{\prime}$ be as in Proposition 23.6. Given $\Gamma$, we can take $\mathfrak{c}$ (employed in the definition of $\Gamma^{\prime}$ ) so that $\Gamma^{\prime} \subset \Gamma$; changing $\mathfrak{c}$ for its suitable multiple, we may also assume that $X \neq \varnothing$, by virtue of [S97, Lemma 11.14 (3)]. By Lemma 20.12 (2), $\mathcal{S}_{k}^{r}\left(C^{r}\right)$ is spanned by eigenforms. Therefore the desired meromorphic continuation follows from Proposition 23.6 and Theorem 23.9. To obtain $E(z, s ; \mathbf{f}, \chi, C)$ from $\mathcal{F}(z, s ; \mathbf{f}, \chi, C)$ whose analytic properties are given in Theorem 23.9, we have to divide the latter by $\mathcal{Z}(u s, \mathbf{f}, \chi)$ and some $L$-functions. Employing Theorem 20.13 and the standard fact on the nonvanishing of $L$-functions, we obtain the last assertion.

For $0<r \in \mathbf{Z}$ and a weight $k$ we put

$$
\Lambda(r, k)=\left\{\begin{array}{lc}
\{x \in \mathbf{R} \mid x \geq 2\} & \text { (Case SP: } \left.r=1 ; k=[k] \text { or } \operatorname{Max}_{v \in \mathbf{a}} k_{v}>3 / 2\right)  \tag{23.30}\\
\{x \in \mathbf{R} \mid x>(3 r / 2)+1\} & \text { (Case SP: all other cases) } \\
\{x \in \mathbf{R} \mid x>3 r\} & \text { (Case UT). }
\end{array}\right.
$$

23.11. Theorem. Suppose that $n>r>0$; put $\lambda=(n+r+1) / 2$ in Case $S P$ and $\lambda=n+r$ in Case UT. Let $k$ be a weight; put $m=k$ in Case SP and $m=\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$ in Case UT. Let $\mathbf{f}$ and $\chi$ be as above with $\kappa=0$ in (16.24a); let $g \in \mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}\right)$ with an arbitrary congruence subgroup $\Gamma$ of $G^{n}$. In Case UT let $\chi_{1}$ denote the restriction of $\chi$ to $F_{\mathbf{A}}^{\times}$. Define $E_{q}$ as in §23.4; define also $\mathcal{F}_{q}$ as in Theorem 23.9 when $\mathbf{f}$ is a Hecke eigenform. Let $\mu \in 2^{-1} \mathbf{Z}$ in Case SP and $\mu \in \mathbf{Z}$ in Case UT.
(I) If $\lambda \leq \mu \leq m_{v}$ and $\mu-m_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$, and $\mu \in \Lambda(r, k)$, then $E_{k}^{n, r}(z, \mu / 2 ; g, \Gamma)$ and $E_{q}(z, \mu / 2 ; \mathbf{f}, \chi, C)$ belong to $\mathcal{N}_{k}^{t}$ except when $F=\mathbf{Q}$ and $\mu=(n+r+2) / 2$ in Case SP, where

$$
t= \begin{cases}(n+r)(m-\mu+2) / 2 & \text { if } \mu=\lambda+1 \text { and } F=\mathbf{Q} \\ (n+r)(m-\mu \mathbf{a}) / 2 & \text { otherwise }\end{cases}
$$

(II) If $\mu$ is as in (I), then $E_{q}(z, \mu / 2 ; \mathbf{f}, \chi, C)$ belongs to $\mathcal{N}_{k}^{t}$ with $t=(n+$ $r)(m-\mu \mathbf{a}) / 2$ except in the following two cases:
(A) Case $\mathrm{SP}, F=\mathbf{Q}, 2 \mu \in\{n+r+2, n+r+3\}$, and $\chi^{2}=1$,
(B) Case UT, $F=\mathbf{Q}, \mu=n+r+1$, and $\chi_{1}=\theta^{\mu}$.
(If $\mu=\lambda+1$, then the statement of (I) is applicable regardless of the nature of $\chi$.)
(III) Suppose that $2 \lambda-m_{v} \leq \mu \leq m_{v}$ and $|\mu-\lambda|+\lambda-m_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$. Then $\mathcal{F}_{q}(z, \mu / 2 ; \mathbf{f}, \chi, C)$ belongs to $\mathcal{N}_{k}^{t}$, where

$$
t=\left\{\begin{array}{l}
(n+r)(k-\mu+2) / 2 \quad\left(\text { Case SP, } \mu=\lambda+1, F=\mathbf{Q}, \text { and } \chi^{2}=1\right) \\
\left.(n+r)(m-\mu+2) / 2 \quad \text { Case UT, } \mu=\lambda+1, F=\mathbf{Q}, \text { and } \chi_{1}=\theta^{\mu}\right) \\
(n+r)\{m-|\mu-\lambda| \mathbf{a}-\lambda \mathbf{a}\} / 2 \text { otherwise }
\end{array}\right.
$$

except in the following four cases:
(C) Case SP, $\mu=0, \mathfrak{c}=\mathfrak{g}$, and $\chi=1$;
(D) Case SP, $0<\mu \leq(n+r) / 2, \mathfrak{c}=\mathfrak{g}$, and $\chi^{2}=1$;
(E) Case SP, $\mu=(n+r+2) / 2, F=\mathbf{Q}$, and $\chi^{2}=1$;
(F) Case UT, $0 \leq \mu<n+r, \mathfrak{c}=\mathfrak{g}$, and $\chi_{1}=\theta^{\mu}$.
23.12. Theorem. Let $n, r, \lambda, \mathbf{f}, \chi$, and $\chi_{1}$ be as in Theorem 23.11 and let $k$ be a weight. Suppose $k=\mu \mathbf{a}$ in Case SP and $\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}=\mu \mathbf{a}$ in Case UT with $\mu$ such that $0<\mu<\lambda$; put $s_{0}=\lambda-(\mu / 2)$. Then $\mathcal{F}_{q}(z, s ; \mathbf{f}, \chi, C)$ has at most a simple pole at $s_{0}$, which occurs only when $\chi^{2}=1$ in Case SP and $\chi_{1}=\theta^{\mu}$ in Case UT. Moreover, the residue is an element of $\mathcal{M}_{k}$.

These two theorems will be proven in §25.7. We can naturally ask whether the nearly holomorphic functions of Theorem 23.11 are arithmetic up to a well-defined constant. That is indeed so in most cases as will be shown in Theorems 27.16 and 28.9. We end this section by proving three lemmas concerning the rationality of certain automorphic forms and Hecke eigenvalues.
23.13. Lemma. Put $\mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}, D\right)=\mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}\right) \cap \mathcal{M}_{k}^{r}(D)$ for any subfield $D$ of $\mathbf{C}$. Let $\Phi$ be the Galois closure of $K$ over $\mathbf{Q}$ in $\mathbf{C}$. Then $\mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}\right)=$ $\mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}, D\right) \otimes_{D} \mathbf{C}$ if $D \supset \Phi \mathbf{Q}_{\mathbf{a b}}$. In particular, suppose $k=\kappa \mathbf{a}$ with $\kappa \in 2^{-1} \mathbf{Z}$ in Case SP and $\kappa \in \mathbf{Z}$ in Case UT. Then $\mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}\right)=\mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}, D\right) \otimes_{D} \mathbf{C}$ if $D \supset \mathbf{Q}_{\mathrm{ab}}$ in Case $S P$ and $D \supset K^{\prime} \mathbf{Q}_{\mathrm{ab}}$ in Case UT, where $K^{\prime}$ is the reflex field defined for $(K, \tau)$ of §3.5 as in §1.12.

Proof. The notation being as in (23.10), let $\varepsilon=\pi_{r}(\gamma)$ with $\gamma \in \Gamma \cap \mathcal{P}^{n, r}$; then $\lambda_{r}(\gamma) \in \mathfrak{r}^{\times}$and $|\varphi(\varepsilon)|=1$. Since $\varphi(\varepsilon)^{2}=\lambda_{r}(\gamma)^{\ell}{\overline{\lambda_{r}(\gamma)}}^{-\ell}$, we see that $\varphi(\varepsilon)^{2}$ is a unit contained in $\Phi$, and $\left|\varphi(\varepsilon)^{2 \sigma}\right|=1$ for every isomorphic embedding $\sigma$ of $\Phi$ into C. Therefore $\varphi(\varepsilon)$ is a root of unity whose square belongs to $\Phi$. Thus $\varphi$ is of finite order. Suppose $D \supset \Phi \mathbf{Q}_{\mathrm{ab}} ;$ let $\Delta=\pi_{r}\left(\Gamma \cap \mathcal{P}^{n, r}\right)$ and $f \in \mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}\right)$. By Theorem 10.8 (2), we can put $f=\sum_{a \in A} a g_{a}$ with a finite subset $A$ of $\mathbf{C}$ and $g_{a} \in \mathcal{S}_{k}^{r}(D)$. We may assume that $A$ is linearly independent over $D$. Then for $\varepsilon \in \Delta$ we have $\sum_{a \in A} a \varphi(\varepsilon) g_{a}=\varphi(\varepsilon) f=f\left\|\varepsilon=\sum_{a \in A} a g_{a}\right\| \varepsilon$. Since $\varphi(\varepsilon) \in \mathbf{Q}_{\mathrm{ab}}$ and $g_{a} \| \varepsilon \in \mathcal{S}_{k}^{r}(D)$ by Theorem 9.13 (3), we have $\varphi(\varepsilon) g_{a}=g_{a} \| \varepsilon$, that is, $g_{a} \in \mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}, D\right)$. This proves the first assertion. Suppose $k=\kappa \mathbf{a}$ with $\kappa$ given as above. Let $D^{\prime}=\mathbf{Q}_{\mathrm{ab}}$ in Case SP and $D^{\prime}=K^{\prime} \mathbf{Q}_{\mathrm{ab}}$ in Case UT. Given $f \in \mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}, \overline{\mathbf{Q}}\right)$, take a finite Galois extension $E$ of $D^{\prime}$ so that $f$ is $E$-rational and put $G=\operatorname{Gal}\left(E / D^{\prime}\right)$. Take $\gamma$ and $\varepsilon$ as above. Then for $\sigma \in G$ we have $\varphi(\varepsilon) f^{\sigma}=(\varphi(\varepsilon) f)^{\sigma}=(f \| \varepsilon)^{\sigma}=$
$f^{\sigma} \| \varepsilon$ by Lemma 10.10, and hence $f^{\sigma} \in \mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}, E\right)$. Therefore $\sum_{\sigma \in G}(b f)^{\sigma} \in$ $\mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}, D^{\prime}\right)$ for every $b \in E$. This shows that $f$ is an $E$-linear combination of elements of $\mathcal{S}_{k}^{r}\left(\Gamma, \mathcal{P}^{n, r}, D^{\prime}\right)$. This combined with the first assertion proves the second assertion.
23.14. Lemma. Let $\mathbf{f} \in \mathcal{M}_{k}(C, \psi)$ as in §20.1 with $\psi$ of finite order, and let $\sigma \in \operatorname{Aut}(\mathbf{C}) ;$ put $\mu(\tau, q ; \mathbf{f})=\left|\operatorname{det}(q)_{\mathbf{h}}\right|_{\mathbf{A}}^{-\nu} c(\tau, q ; \mathbf{f})$, where $\nu=0$ or $1 / 2$ according as $k$ is integral or half-integral. Then there exists an element of $\mathcal{M}_{k^{\sigma}}\left(C, \psi^{\sigma}\right)$, which we write $\mathbf{f}^{\sigma}$ and which is uniquely determined by the property $\mu\left(\tau, q ; \mathbf{f}^{\sigma}\right)=\mu(\tau, q ; \mathbf{f})^{\sigma}$ for every $(\tau, q)$, where $k^{\sigma}$ is defined as in Theorem 10.4 (5) or Theorem 10.7 (5). Moreover, if the conditions of Lemma 20.8 are satisfied and $\mathbf{f} \mid T(\mathfrak{a})=\lambda(\mathfrak{a}) \mathbf{f}$ as in §20.6 or §21.3, then $\mathbf{f}^{\sigma} \mid T(\mathfrak{a})=\lambda_{\sigma}(\mathfrak{a}) \mathbf{f}^{\sigma}$, where $\lambda_{\sigma}(\mathfrak{a})=\lambda(\mathfrak{a})^{\sigma}\left(N(\mathfrak{a})^{\nu}\right)^{\sigma} / N(\mathfrak{a})^{\nu}$.

Proof. Let $f_{p}$ be the $p$-component of $\mathbf{f}$ with $p=\operatorname{diag}[q, \widetilde{q}], q \in G L_{n}(K)_{\mathbf{h}}$. By Theorem 10.4 (2), Theorem 10.7 (2), and the remark after the proof of Theorem 10.4, we have $\left(f_{p}\right)^{\sigma} \in \mathcal{M}_{k^{\sigma}}\left(\Gamma^{p}, \psi^{\sigma}\right)$, where $\Gamma^{p}=G \cap p C p^{-1}$ and we take $p=1$ if $k \notin$ $\mathbf{Z}^{\mathrm{b}}$. (Take $W$ of Theorem 10.4 to be $p C p^{-1}$. Then $\Gamma$ there is $\Gamma^{p}$.) We first consider the case of integral $k$. By [S97, Lemma $8.8(3)]$ there is a finite subset $Q$ of $G L_{n}(K)_{\mathbf{h}}$ such that $G_{\mathbf{A}}=\bigsqcup_{q \in Q} G \operatorname{diag}[q, \widehat{q}] C$; in Case SP strong approximation allows us to take $Q=\{1\}$. Given $\mathbf{f}$, we define $\mathbf{f}^{\sigma} \in \mathcal{M}_{k^{\sigma}}\left(C, \psi^{\sigma}\right)$ so that its $p$-component is $\left(f_{p}\right)^{\sigma}$ for every $p=\operatorname{diag}[q, \widehat{q}]$ with $q \in Q$. Our task is to show that $c\left(\tau, q ; \mathbf{f}^{\sigma}\right)=$ $c(\tau, q ; \mathbf{f})^{\sigma}$ for every $(\tau, q)$. Since $\mathcal{M}_{k}$ is spanned by $\mathcal{M}_{k}(\overline{\mathbf{Q}})$, it is sufficient to prove the case where $f_{p} \in \mathcal{M}_{k}(\overline{\mathbf{Q}})$ for every such $p$. Given $r \in G L_{n}(K)_{\mathbf{h}}$, we can find $\beta \in G$ and $p=\operatorname{diag}[q, \widehat{q}]$ with $q \in Q$ so that $\operatorname{diag}[r, \widehat{r}] \in \beta p C$. Then we have (20.9e). Take $s \in \mathbf{Z}_{\mathbf{h}}^{\times}$so that $[s, \mathbf{Q}]=\sigma$ on $\mathbf{Q}_{\mathrm{ab}}$, and put $x=\operatorname{diag}\left[1_{n}, s^{-1} 1_{n}\right]$ and $\alpha=\beta^{-1}$. Since $\operatorname{det}\left(x^{-1} \alpha x \alpha^{-1}\right)=1$, by strong approximation in $G_{1}$ we have $x^{-1} \alpha x \alpha^{-1} \in U^{N} G_{1}$ with $U^{N}$ of (8.5) for any $N$. Thus we can put $x^{-1} \alpha x=u \varepsilon$ with $\varepsilon \in G_{1} \alpha$ and $u \in U_{N}$. Then $u \in G_{\mathbf{A}}$. Take $N$ so that both $f_{p}$ and $f_{p} \| \alpha$ belong to $\mathcal{M}_{k}\left(\Gamma^{N}, \overline{\mathbf{Q}}\right)$ with $\Gamma^{N}$ of (7.6). Then $\left(f_{p} \| \alpha\right)^{\sigma}=\left(f_{p}\right)^{(\alpha x, \sigma)}$ and $\left(f_{p}\right)^{\sigma}=\left(f_{p}\right)^{(x u, \sigma)}$ by Theorem 10.2 (8). Thus $\left(f_{p} \| \alpha\right)^{\sigma}=\left(f_{p}\right)^{(\alpha x, \sigma)}=\left(f_{p}\right)^{(x u \varepsilon, \sigma)}=\left(f_{p}\right)^{\sigma} \| \varepsilon$. Put $\gamma=\varepsilon^{-1}$. Since $x C x^{-1}=C$ and $x$ commutes with $p$ and $\operatorname{diag}[r, \widehat{r}]$, we have $\operatorname{diag}[r, \hat{r}] \in x^{-1} \beta x p C=\gamma u^{-1} p C$. Now we can choose sufficiently large $N$ so that $U^{N} \cap G_{\mathbf{A}} \subset p C p^{-1}$. Then $\operatorname{diag}[r, \widehat{r}] \in \gamma p C$, so that we can write formula (20.9e) with ( $\mathbf{f}^{\sigma}, \psi^{\sigma}, \gamma$ ) as ( $\mathbf{f}, \psi, \beta$ ). Applying $\sigma$ to (20.9e), we obtain $\left(f_{p}\right)^{\sigma} \| \gamma^{-1}$ on the left-hand side. Comparing the two equalities, we find that $c\left(\tau, q ; \mathbf{f}^{\sigma}\right)=c(\tau, q ; \mathbf{f})^{\sigma}$, since $u \in U_{N}$ and so we have $\psi_{\mathrm{c}}\left(\operatorname{det}\left(a_{\beta p}^{-1} a_{\gamma p}\right)\right)=1$ for a sufficiently large $N$.

Next let us consider the case of half-integral $k$. We define $\mathbf{f}^{\sigma} \in \mathcal{M}_{k}\left(C, \psi^{\sigma}\right)$ so that its 1-component is $\left(f_{1}\right)^{\sigma}$. Given $r \in G L_{n}(K)_{\mathbf{h}}$, take $\beta \in G$ so that $\operatorname{diag}[r, \widehat{r}] \in$ $\beta C$, and take $\alpha, \varepsilon$, and $\gamma$ as above. Then both $\beta$ and $\gamma$ belong to $\mathfrak{M}$, and so $h_{\beta}$ and $h_{\gamma}$ are meaningful. Put $\theta(z)=\theta(0, z ; \ell)$ with the notation of (A2.23), where we take $\ell$ to be the characteristic function of $\prod_{v \in \mathrm{~h}}\left(\mathfrak{g}_{n}^{1}\right)_{v}$. Then $\theta^{-1} f_{1} \in \mathcal{A}_{[k]}(\overline{\mathbf{Q}})$, and the above reasoning shows that $\left(\left(\theta^{-1} f_{1}\right) \|_{[k]} \beta^{-1}\right)^{\sigma}=\left(\theta^{-1} f_{1}\right)^{\sigma} \|_{\left[k^{\sigma}\right]} \gamma^{-1}$. By Proposition A2.5, $h_{\beta}\left(\beta^{-1} z\right) \theta\left(\beta^{-1} z\right)=\theta\left(0, z ;{ }^{\beta} \ell\right)$, and by Theorem A2.4 (6), (8) and (A2.3a) we have $\left({ }^{\beta} \ell\right)(y)=e \ell(y r)$, where $e=|\operatorname{det}(r)|_{\mathbf{A}}^{1 / 2}$. The same holds with $\gamma$ in place of $\beta$, so that $\left[e^{-1} h_{\beta}\left(\beta^{-1} z\right) \theta\left(\beta^{-1} z\right)\right]^{\sigma}=\theta\left(0, z ; e^{-1} \cdot{ }^{\beta} \ell\right)^{\sigma}=\theta\left(0, z ; e^{-1} \cdot \beta \ell\right)=$ $e^{-1} h_{\gamma}\left(\gamma^{-1} z\right) \theta\left(\gamma^{-1} z\right)$, since $e^{-1} \cdot{ }^{\beta} \ell$ is $\mathbf{Q}$-valued. Therefore we obtain

$$
\left[e^{-1} j_{\beta}^{k}\left(\beta^{-1} z\right) f_{1}\left(\beta^{-1} z\right)\right]^{\sigma}=e^{-1} j_{\gamma}^{k^{\top}}\left(\gamma^{-1} z\right) f_{1}^{\sigma}\left(\gamma^{-1} z\right)
$$

Employing (20.9e), we obtain $\left[e^{-1} c(\tau, r ; \mathbf{f})\right]^{\sigma}=e^{-1} c\left(\tau, r ; \mathbf{f}^{\sigma}\right)$ as expected.

To prove the assertion concerning $\lambda(\mathfrak{a})$, substitute $N(\mathfrak{a})^{\nu}[\mathfrak{a}]$ for [a] in (20.18), and apply $\sigma$ to the coefficients of the formal Dirichlet series; see Lemma 20.5 and Theorem 21.4 for the explicit forms of $A(\tau, L)$. Comparing the coefficients of [a], we obtain the last equality of our proposition.
23.15. Lemma. The notation and assumption being the same as in Lemma 23.14, suppose that $\mathbf{f} \mid T(\mathfrak{a})=\lambda(\mathfrak{a}) \mathbf{f}$ for every $\mathfrak{a}$. Let $W$ be the field generated by the $\lambda(\mathfrak{a})$ over $\mathbf{Q}$ for all $\mathfrak{a}$ prime to $\mathfrak{c}$. Then $W$ is totally real or a CM-field, and the latter can occur only in Case UT.

Proof. Denote complex conjugation in $\mathbf{C}$ by $\rho$. If $T(\mathfrak{a})=\sum_{\tau} C \tau C$ and $\tau_{v}=1$ for every $v \mid \mathbf{c}$, then we easily see that $T\left(\mathfrak{a}^{\rho}\right)=\sum_{\tau} C \tau^{-1} C$. Therefore from [S97, Proposition 11.7] and [S95b, Lemma 4.5] we can easily derive that $\lambda(\mathfrak{a})^{\rho}=\lambda\left(\mathfrak{a}^{\rho}\right)$ if $\mathfrak{a}$ is prime to $\mathfrak{c}$. In Case SP this means that $\lambda(\mathfrak{a})$ is real, and hence, by Lemma 23.14, $\lambda_{\sigma}(\mathfrak{a})$ is real for every $\sigma \in \operatorname{Aut}(\mathbf{C})$. Since $\left(N(\mathfrak{a})^{\nu}\right)^{\sigma}= \pm N(\mathfrak{a})^{\nu}$, we see that $\lambda(\mathfrak{a})^{\sigma}$ is real for every $\sigma$, so that $W$ is totally real. In Case UT, we have $\lambda(\mathfrak{a})^{\sigma \rho}=\lambda_{\sigma}(\mathfrak{a})^{\rho}=\lambda_{\sigma}\left(\mathfrak{a}^{\rho}\right)=\lambda\left(\mathfrak{a}^{\rho}\right)^{\sigma}=\lambda(\mathfrak{a})^{\rho \sigma}$, so that $\sigma \rho=\rho \sigma$ on $W$. Thus $W$ is totally real or a CM-field in Case UT.

## 24. Pullback of Eisenstein series

We fix two integers $r$ and $n$ as in $\S 23.1$ and put $t=n-r$ and $N=n+r$. Taking ( $n, r, N$ ) in place of ( $r, n-r, n$ ) in (23.5), for $(\beta, \gamma) \in G^{n} \times G^{r}$ we can define $\beta \times \gamma$ as an element of $G^{N}$ and view $G^{n} \times G^{r}$ as a subgroup of $G^{N}$.
24.1. Lemma. We have $G^{N}=\bigsqcup_{\nu=0}^{r} P^{N} \tau_{\nu}\left(G^{n} \times G^{r}\right)$ with $\tau_{\nu}$ given by

$$
\tau_{\nu}=\left[\begin{array}{cc}
1_{N} & 0  \tag{24.1}\\
f_{\nu} & 1_{N}
\end{array}\right], \quad f_{\nu}=\left[\begin{array}{cc}
0 & g_{\nu} \\
t g_{\nu} & 0
\end{array}\right], \quad g_{\nu}=\left[\begin{array}{cc}
1_{\nu} & 0 \\
0 & 0
\end{array}\right] \in K_{r}^{n} .
$$

Moreover $P^{N} \tau_{\nu}\left(G^{n} \times G^{r}\right)=\bigsqcup_{\xi, \beta, \gamma} P^{N} \tau_{\nu}\left(\left(\xi \times 1_{2 n-2 \nu}\right) \beta \times \gamma\right)$, where $\xi$ runs over $G^{\nu}, \beta$ over $P^{n, \nu} \backslash G^{n}$, and $\gamma$ over $P^{r, \nu} \backslash G^{r}$. Furthermore, $\left(\kappa_{\nu} \pi_{\nu}(\gamma) \kappa_{\nu} \times 1_{2 n-2 \nu}\right) \times$ $\gamma \in \tau_{\nu}^{-1} P^{N} \tau_{\nu}$ for every $\gamma \in P^{r, \nu}$, where $\kappa_{\nu}=\left[\begin{array}{cc}0 & 1_{\nu} \\ 1_{\nu} & 0\end{array}\right]$.

Proof. This is essentially included in [S97, Propositions 2.4 and 2.7, and Lemma 2.6]. Indeed, put $\varphi=\eta_{r}$,

$$
\begin{gather*}
\omega=\left[\begin{array}{cc}
\psi & 0 \\
0 & -\varphi
\end{array}\right], \quad \psi=\left[\begin{array}{ccc}
0 & 0 & -1_{t} \\
0 & \eta_{r} & 0 \\
1_{t} & 0 & 0
\end{array}\right], \quad R=\operatorname{diag}\left[1_{n},\left[\begin{array}{cc}
0 & 1_{r} \\
1_{n} & 0
\end{array}\right], 1_{r}\right],  \tag{24.2}\\
T=\left[\begin{array}{ccc}
0 & 1_{r} & 0 \\
1_{t} & 0 & 0 \\
0 & 0 & 1_{n}
\end{array}\right], \quad S=\left[\begin{array}{cccc}
1_{t} & 0 & 0 & 0 \\
0 & 1_{2 r} & 0 & -\lambda \\
0 & 0 & 1_{t} & 0 \\
0 & -1_{2 r} & 0 & \lambda^{*}
\end{array}\right], \quad \lambda=\left[\begin{array}{cc}
0_{r} & 1_{r} \\
0_{r} & 0_{r}
\end{array}\right] \tag{24.3}
\end{gather*}
$$

Denote the group $U(\varphi)$ of (1.7) by $G^{\varphi}=G(\varphi)$ in conformity with the notation of [S97]. We have $G^{n}=G\left(\eta_{n}\right)$, for example. Then $\eta_{r}=\lambda^{*}-\lambda$ and $S \eta_{N} S^{*}=\omega$, and so $S^{-1} G^{\omega} S=G^{N} ; \psi={ }^{t} T \eta_{n} T$, and so $G^{\psi}=T^{-1} G^{n} T$; also $R \cdot \operatorname{diag}[\alpha, \beta] R^{-1}=$ $\alpha \times \beta$ for $\alpha \in G^{n}$ and $\beta \in G^{r}$. Now in [S97, Proposition 2.4] we showed that $P_{U}^{\omega} \backslash G^{\omega} /\left[G^{\psi} \times G^{\varphi}\right]$ for a certain parabolic subgroup $P_{U}^{\omega}$ of $G^{\omega}$ has exactly $r+$ 1 orbits, say $X_{\nu}$ for $0 \leq \nu \leq r$, and gave an explicit set of representatives for $P_{U}^{\omega} \backslash X_{\nu}$. Observing that $P_{U}^{\omega}=S P^{N} S^{-1}$ (cf. [S97, (21.1.8)]) and employing the above isomorphisms among the groups involved, we see that $P^{N} \backslash G^{N} /\left[G^{n} \times G^{r}\right]$
has exactly $r+1$ orbits. To find good representatives, for $x \in G^{N}$ put $\left[\begin{array}{cc}c_{x} & d_{x}\end{array}\right]=$ $\left[\begin{array}{llll}a & b & a^{\prime} & b^{\prime}\end{array}\right]$ with $a, a^{\prime} \in K_{n}^{N}$ and $b, b^{\prime} \in K_{r}^{N}$. Clearly the rank of $\left[\begin{array}{ll}a & a^{\prime}\end{array}\right]$ depends only on $P^{N} x\left(G^{n} \times G^{r}\right)$. Since this rank is $n+\nu$ if $x=\tau_{\nu}$, we obtain the first asserion of our proposition. Then the second and third assertions can easily be verified by translating (or by modifying) the proof of [S97, Lemma 2.6 and Proposition 2.7]. For a more direct proof in Case SP, see [S95a, pp.556-557].

We now consider $C^{r}$ and $C^{n, r}$ of $\S 23.3$. We put $D^{n}=D^{n}\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]$ as we did there; hereafter we assume that $\mathfrak{e}=\mathfrak{c}$. We define a subgroup $C^{\prime}$ of $D^{n}$ and an element $\sigma$ of $G_{\mathrm{h}}^{N}$ by

$$
\begin{gather*}
C^{\prime}=\left\{x \in D^{n} \mid d_{3}(x) \prec \mathfrak{r c}\right\},  \tag{24.4}\\
\sigma_{v}= \begin{cases}\operatorname{diag}\left[1_{n}, \theta_{v}^{-1} 1_{r}, 1_{n}, \theta_{v} 1_{r}\right] \tau_{r} & \text { if } v \mid \mathfrak{c}, \\
\operatorname{diag}\left[1_{n}, \theta_{v}^{-1} 1_{r}, 1_{n}, \theta_{v} 1_{r}\right] & \text { if } v \nmid \mathfrak{c},\end{cases} \tag{24.5}
\end{gather*}
$$

where $\theta$ is an element of $F_{\mathbf{h}}^{\times}$such that $\theta \mathfrak{g}=\mathfrak{b}$. To see that $C^{\prime}$ is indeed a subgroup, take the homomorphism $x \mapsto\left(\left(d_{x}\right)_{v}\right)_{v \mid \mathfrak{c}}$ of $D^{n}$ into $\prod_{v \mid \mathfrak{c}} G L_{n}\left(\mathfrak{r}_{v} / \mathfrak{r}_{v} \mathfrak{c}_{v}\right)$ noted in (1.18). Then $C^{\prime}$ is the inverse image of the subgroup of the latter group defined by the vanishing of the lower left $(t \times r)$-block. Notice also that $C^{\prime}$ can be defined by the condition $a_{2}(x) \prec \mathfrak{r c}$ instead of $d_{3}(x) \prec \mathfrak{r c}$.
24.2. Lemma. (1) $P_{\mathbf{A}}^{n, r} C^{n, r}=P_{\mathbf{A}}^{n, r} C^{\prime}$.
(2) Let $q$ be any fixed element of $G_{\mathbf{h}}^{n}$ as in (23.25). Then

$$
\left(P^{N} \tau_{r}\left(G^{n} \times G^{r}\right)\right) \cap P_{\mathbf{A}}^{N} D^{N} \sigma=\bigsqcup_{\xi \in \Xi, \beta \in \mathcal{R}_{q}} P^{N} \tau_{r}\left(\left(\xi \times 1_{2 t}\right) \beta \times 1_{2 r}\right)
$$

where $\mathcal{R}_{q}$ is the subset of $\bigsqcup_{\zeta} G^{n} \cap\left(\zeta q C^{n, r} q^{-1}\right)$ of $\S 23.4$, and $\Xi=G^{r} \cap \mathfrak{X}$ with $\mathfrak{X}$ of (19.2c) defined with $\mathfrak{e}=\mathfrak{c}$ and $r$ in place of $n$ there; we take $\mathcal{R}_{q}=\{1\}$ if $r=n$.

Proof. To prove (1), by mens of the map $x \mapsto \varepsilon x \varepsilon^{-1}$ with $\varepsilon=\operatorname{diag}\left[1_{n}, \theta 1_{n}\right]$, we may assume that $\mathfrak{b}=\mathfrak{g}$. Given $x \in C^{\prime}$, put $y=\pi_{r}(x)$. Then we easily see that $y \eta_{r} y^{*}-\eta_{r}=z-z^{*}$ with a matrix $z \prec \mathfrak{r c}$. Fix a prime $v \mid \mathfrak{c}$, assume $K \neq F$, and put $f=e-e^{\rho}$ with $e$ such that $\mathfrak{r}_{v}=\mathfrak{g}_{v}[e]$. Then $\left(z-z^{*}\right)_{v} \in f \mathfrak{c}_{v} \widetilde{S}_{v}$ with $\widetilde{S}_{v}$ of (16.1c). Then successive approximation produces an element $k_{v}$ of $G_{v}^{r} \cap G L_{2 r}\left(\mathfrak{r}_{v}\right)$ such that $k_{v}^{-1} y_{v}-1 \prec \mathfrak{r}_{v} \mathfrak{c}_{v}$, as proved in [S97, Lemma 17.2 (2)]. (Take $(\delta, \varphi)$ there to be ( $f, f^{-1} \eta_{r}$ ) here.) Let $p=\left(p_{v}\right)$ with $p_{v}=k_{v} \times 1_{2 t}$ for $v \mid c$ and $p_{v}=1$ for all other $v$ 's. Then $p \in P_{\mathrm{h}}^{n, r}$ and $p^{-1} x \in C^{n, r}$, and hence $x \in P_{\mathrm{A}}^{n, r} C^{n, r}$, which proves (1) in Case UT, since $C^{n, r} \subset C^{\prime}$. Case SP, in which $f$ is unnecessary, can be proved in a similar and simpler way.

Next, to prove (2), we first assume $q=1$. Let $\alpha=\left(\xi \times 1_{2 t}\right) \beta$ with $\xi \in G^{r}$ and $\beta \in P^{n, r} \backslash G^{n}$. By Lemma 24.1, $P^{N} \tau_{r}\left(G^{n} \times G^{r}\right)$ is a disjoint union of $P^{N} \tau_{r}\left(\alpha \times 1_{2 r}\right)$ with such $\alpha$ 's. Thus our task is to determine $P^{N} \tau_{r}\left(\alpha \times 1_{2 r}\right)$ contained in $P_{\mathbf{A}}^{N} D^{N} \sigma$. Put $\omega=\sigma\left(\alpha \times 1_{2 r}\right) \sigma^{-1}$. Since $G_{v}^{N}=P_{v}^{N} D_{v}^{N}$ for $v \nmid \mathfrak{c}$, and $\left(\tau_{r} \sigma^{-1}\right)_{v} \in P_{v}^{N}$ for $v \mid \mathbf{c}$, we only have to find those $\alpha$ such that $\omega \in P_{\mathbf{A}}^{N} D^{N}$. Clearly $\omega_{v}=\left(\alpha \times 1_{2 r}\right)_{v}$ for $v \nmid c$. Fix a prime $v \mid c$ and write $\alpha$ in the form (23.1). Then

$$
\omega_{v}=\left[\begin{array}{cccccc}
a_{1} & a_{2} & -\theta b_{1} & b_{1} & b_{2} & 0  \tag{24.6}\\
a_{3} & a_{4} & -\theta b_{3} & b_{3} & b_{4} & 0 \\
0 & 0 & 1 & 0 & 0 & 0 \\
c_{1} & c_{2} & \theta\left(1-d_{1}\right) & d_{1} & d_{2} & 0 \\
c_{3} & c_{4} & -\theta d_{3} & d_{3} & d_{4} & 0 \\
\theta\left(a_{1}-1\right) & \theta a_{2} & -\theta^{2} b_{1} & \theta b_{1} & \theta b_{2} & 1
\end{array}\right]_{v} .
$$

By Lemma 1.9, $\omega \in P_{\mathbf{A}}^{N} D^{N}$ if and only if $\left(d_{\omega}\right)_{v} \in G L_{N}\left(K_{v}\right)$ and $\left(d_{\omega}^{-1} c_{\omega}\right)_{v} \prec(\mathfrak{r b c})_{v}$ for every $v \mid \mathfrak{c}$, which is so only if $d_{\alpha} \in G L_{n}\left(K_{v}\right)$ and $d_{\alpha}^{-1} c_{\alpha} \prec(\mathfrak{r b c})_{v}$ for every $v \mid \mathfrak{c}$, as can easily be seen from (24.6). Then $\alpha \in P_{\mathbf{A}}^{n} D^{n} \subset P_{\mathbf{A}}^{n, r} D^{n}$ by Lemmas 1.9 and 23.7 , and $\beta \in P_{\mathbf{A}}^{n, r} D^{n}$, since $\xi \times 1 \in P^{n, r}$. Thus we may restrict $\beta$ to $G^{n} \cap P_{\mathbf{A}}^{n, r} D^{n}$. As explained in §23.4, we may assume that $\beta=\zeta w$ with $w \in D^{n}$ and $\zeta$ as in (23.27). If $r=n$, we can take $\beta=1$, and so can take $\zeta=w=1$. Assuming $n>r$, we have $\xi \times 1 \in P_{\mathbf{A}}^{n} D^{n}$ since $(\xi \times 1) \zeta w=\alpha \in P_{\mathbf{A}}^{n} D^{n}, \zeta_{v}=1$ for $v \mid \mathfrak{c}$, and $P_{v}^{n} D_{v}^{n} \neq G_{v}^{n}$ only if $v \mid \mathfrak{c}$. Consequently $d_{\xi} \in G L_{r}\left(K_{v}\right)$ and $d_{\xi}^{-1} c_{\xi} \prec(\mathfrak{r b c})_{v}$ for every $v \mid \mathfrak{c}$ by Lemma 1.9. Write $\zeta w$ in the form (23.1) and put $p_{i}=a(\zeta w)_{i}, q_{i}=$ $b(\zeta w)_{i}, r_{i}=c(\zeta w)_{i}, s_{i}=d(\zeta w)_{i}$; put $e=d_{\xi}^{-1} c_{\xi}$. Then
$\left[\begin{array}{cc}c_{\alpha} & \left.d_{\alpha}\right]=\operatorname{diag}\left[d_{\xi}, 1_{2 t}\right]\left[\begin{array}{ll}r^{\prime} & \left.s^{\prime}\right]\end{array} \quad \text { with } \quad r^{\prime}=r+\operatorname{diag}[e, 0] p \quad \text { and } \quad s^{\prime}=s+\operatorname{diag}[e, 0] q . ~\right.\end{array}\right.$
Compute $c_{\omega}$ and $d_{\omega}$ with these $c_{\alpha}$ and $d_{\alpha}$. Fix $v \mid \boldsymbol{c}$. Recall that $\zeta_{v}=1$ for such a $v$. Since $e_{v} \prec(\mathfrak{r b c})_{v}, q_{v} \prec\left(\mathfrak{r b}^{-1}\right)_{v}$ and $s_{v} \in G L_{r}\left(\mathfrak{r}_{v}\right)$, we see that $s_{v}^{\prime} \in G L_{r}\left(\mathfrak{r}_{v}\right)$. Focusing our attention on the upper right $(n \times r)$-block of $\left(d_{\omega}^{-1} c_{\omega}\right)_{v}$ we see that

$$
s_{v}^{\prime-1}\left[\begin{array}{c}
\theta\left(d_{\xi}^{-1}-s_{1}^{\prime}\right)  \tag{*}\\
-\theta s_{3}
\end{array}\right]_{v} \prec(\mathfrak{r b c})_{v} .
$$

Thus $\left(s_{3}\right)_{v} \prec(\mathfrak{r c})_{v}$, and hence $w \in C^{\prime}$, so that $\beta \in \zeta C^{\prime} \subset P_{\mathbf{A}}^{n, r} C^{\prime}=P_{\mathbf{A}}^{n, r} C^{n, r}$ by (1). This means that changing $w$ suitably, we may now assume that $w \in C^{n, r}$. From (24.6) we easily see that $\sigma(w \times 1) \sigma^{-1} \in D^{N}$, and hence $\sigma((\xi \times 1) \zeta \times 1) \sigma^{-1} \in$ $P_{\mathbf{A}}^{N} D^{N}$ which is true also when $n=r$, since $\zeta=w=1$. Thus, for $\beta \in \zeta C^{n, r}$ we have $\sigma((\xi \times 1) \beta \times 1) \sigma^{-1} \in P_{\mathbf{A}}^{N} D^{N}$ if and only if $\left[\sigma((\xi \times 1) \times 1) \sigma^{-1}\right]_{v} \in P_{v}^{N} D_{v}^{N}$ for every $v \mid c$, in which case we can repeat the above computation of $d_{\omega}^{-1} c_{\omega}$ with $w=1$. Then in (*) we have $s_{v}^{\prime}=1$, and hence $d_{\xi}-1 \prec \mathfrak{r}_{v} \mathfrak{c}_{v}$ and $c_{\xi} \prec(\mathfrak{r b c})_{v}$. From (24.6) we see that the lower right ( $r \times r$ )-block of $d_{\omega}^{-1} c_{\omega}$ is $-\theta^{2} b_{\xi} d_{\xi}^{-1}$, and so $b_{\xi} \prec\left(\mathfrak{r b}^{-1} \mathfrak{c}\right)_{v}$. Thus $\xi \in \mathfrak{X}$. Conversely suppose $\xi \in G^{r} \cap \mathfrak{X}$ and $v \mid \mathfrak{c}$. Then $\xi_{v} \times 1 \in C_{v}^{n, r}$, and so $\left[\sigma((\xi \times 1) \times 1) \sigma^{-1}\right]_{v} \in D_{v}^{N}$. This proves (2) when $q=1$. For a more general $q$, we can repeat our argument with $\beta \in \zeta q C^{n, r} q^{-1}$ since $q_{v}=1$ for every $v \mid c$. This completes our proof.
24.3. Lemma. Define $\mathrm{il}_{\mathfrak{b}}$ on $G L_{2 n+2 r}(K)_{\mathbf{A}}$ by (1.19) (or (16.23)) by taking $N$ as both $m$ and $n$ there. Let $\alpha=\tau_{r}\left(\left(\xi \times 1_{2 t}\right) \beta \times 1_{2 r}\right)$ with $\xi \in \Xi$ and $\beta \in \zeta q C^{n, r} q^{-1}$, where $q=g \times \operatorname{diag}[\widehat{\varphi}, \varphi]$ and $\zeta=e \times \operatorname{diag}[\widehat{s}, s]$ as in (23.25) and (23.27); let $x=\alpha(q \times \varepsilon h \varepsilon) \sigma^{-1}$ with $\varepsilon=\left[\begin{array}{ll}0 & 1 \\ 1 & 0\end{array}\right]$ and $h \in G_{\mathrm{h}}^{r}$ such that $h_{v}=1$ for every $v \mid \mathfrak{c}$. Then $\nu_{\mathfrak{b}}\left(h^{-1} \xi e g\right)$ is prime to $\mathfrak{c}$ and

$$
\begin{gather*}
\mathrm{il}_{\mathfrak{b}}(x)=\theta^{-r} \operatorname{det}(h) \operatorname{det}(s \varphi) \nu_{\mathfrak{b}}\left(h^{-1} \xi e g\right)^{-1}  \tag{24.7}\\
{\left[\operatorname{det}\left(d_{x}\right)^{-1} \theta^{-r} \lambda_{r}(\beta)\right]_{v} \equiv 1\left(\bmod \mathfrak{c}_{v} \mathfrak{r}_{v}\right) \text { for every } v \mid \mathfrak{c}} \tag{24.8}
\end{gather*}
$$

where $\nu_{\mathrm{b}}$ is defined by (19.5) and $\lambda_{r}$ is defined by (23.3).
Proof. Put $\beta=\zeta q w q^{-1}$ with $w \in C^{n, r}, h^{\prime}=\varepsilon h \varepsilon$, and $f=\tau_{r}\left((\xi \times 1) \zeta q \times h^{\prime}\right) \sigma^{-1}$. Then $x=f \sigma(w \times 1) \sigma^{-1}$ and $\sigma(w \times 1) \sigma^{-1} \in D^{N}$ as seen in the proof of Lemma 24.2. Therefore $\mathrm{il}_{\mathfrak{b}}(x)=\mathrm{il}_{\mathfrak{b}}(f)$. If $v \mid \mathfrak{c}$, then $\zeta_{v}, q_{v}, h_{v}, e_{v}$, and $g_{v}$ are all identity matrices (of various sizes), and so $f_{v}=\left[\tau_{r} \sigma^{-1} \sigma((\xi \times 1) \times 1) \sigma^{-1}\right]_{v}$. Now $\xi_{v} \times 1 \in$ $C^{n, r}$ since $\xi_{v} \in C^{r}$. Therefore $\operatorname{il}_{\mathfrak{b}}(f)_{v}=\operatorname{il}_{\mathfrak{b}}\left(\tau_{r} \sigma^{-1}\right)_{v}=\theta_{v}^{-r} \mathfrak{r}_{v}$, since

$$
\begin{equation*}
\left(\tau_{r} \sigma^{-1}\right)_{v}=\operatorname{diag}\left[1_{n}, \theta_{v} 1_{r}, 1_{n}, \theta_{v}^{-1} 1_{r}\right] \quad \text { if } \quad v \mid \mathbf{c} . \tag{24.9}
\end{equation*}
$$

Also $\nu_{\mathfrak{b}}\left(\xi_{v}\right)=\mathfrak{r}_{v}$, since $\xi_{v} \in C^{r}$. Thus $\nu_{\mathfrak{b}}\left(h^{-1} \xi e g\right)$ is prime to $c$. Suppose $v \nmid c$. Then, putting $\xi e g=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right]$ and $h^{\prime}=\left[\begin{array}{ll}a^{\prime} & b^{\prime} \\ c^{\prime} & d^{\prime}\end{array}\right]$, by a direct calculation we can show that

$$
\left[\begin{array}{ll}
c_{f} & d_{f}
\end{array}\right]_{v}=\left[\begin{array}{cccccc}
c & 0 & \theta a^{\prime} & d & 0 & \theta^{-1} b^{\prime} \\
0 & 0 & 0 & 0 & s \varphi & 0 \\
a & 0 & \theta c^{\prime} & b & 0 & \theta^{-1} d^{\prime}
\end{array}\right]_{v}
$$

Clearly $\mathrm{il}_{\mathfrak{b}}(f)_{v}=\operatorname{det}(s \varphi)_{v} \mathrm{il}_{\mathfrak{b}}(z)_{v}$ with a matrix $z \in G L_{4 r}\left(K_{v}\right)$ such that

$$
\left[\begin{array}{ll}
c_{z} & d_{z}
\end{array}\right]_{v}=\left[\begin{array}{cccc}
a & \theta c^{\prime} & b & \theta^{-1} d^{\prime} \\
c & \theta a^{\prime} & d & \theta^{-1} b^{\prime}
\end{array}\right]_{v}
$$

Put $p=h \varepsilon, \gamma=h^{-1} \xi e g$, and $A=\operatorname{diag}\left[1_{r},\left[\begin{array}{cc}0 & -\theta^{-1} 1_{r} \\ \theta 1_{r} & 0\end{array}\right], 1_{r}\right]$. Then

$$
p_{v}^{-1}\left[\begin{array}{ll}
c_{z} & d_{z}
\end{array}\right]_{v} A_{v}=\left[\begin{array}{cccc}
c_{\gamma} & \theta d_{\gamma} & -1_{r} & 0 \\
a_{\gamma} & \theta b_{\gamma} & 0 & \theta^{-1} 1_{r}
\end{array}\right]_{v}
$$

Since $A \in D^{2 r}\left[\mathfrak{b}^{-1}, \mathfrak{b}\right]$, by Lemma 1.11 (2) we obtain $\mathrm{il}_{\mathfrak{b}}(z)_{v}=\operatorname{det}(h)_{v} \theta_{v}^{-r} \nu_{0}\left(\gamma^{\prime}\right)_{v}^{-1}$ with $\gamma^{\prime}=\left[\begin{array}{cc}-\theta^{-1} c_{\gamma} & -d_{\gamma} \\ a_{\gamma} & \theta b_{\gamma}\end{array}\right]$. In view of (19.5) we easily see that $\nu_{0}\left(\gamma^{\prime}\right)_{v}=\nu_{b}(\gamma)_{v}$. Combining all these, we obtain (24.7).

Next, let $v \mid \mathfrak{c}$; then $x_{v}=\left(\tau_{r} \sigma^{-1} \omega\right)_{v}$ with $\omega=\sigma((\xi \times 1) \beta \times 1) \sigma^{-1}$. By (24.9), $\operatorname{diag}\left[1_{n}, \theta_{v} 1_{r}\right]\left(d_{x}\right)_{v}=\left(d_{\omega}\right)_{v}$, and $\left(d_{\omega}\right)_{v}$ can be obtained from (24.6) by taking $\alpha$ there to be $(\xi \times 1) \beta$ here. Since $\beta_{v}=w_{v} \in C_{v}^{n, r}$ and $\xi_{v} \in C_{v}^{r}$, we see that

$$
\left(d_{\omega}\right)_{v} \equiv\left[\begin{array}{ccc}
1_{r} & y_{1} & 0 \\
0 & d_{4}(\beta)_{v} & 0 \\
y_{2} & y_{3} & 1_{r}
\end{array}\right] \quad\left(\bmod \mathfrak{r}_{v} \mathfrak{c}_{v}\right)
$$

with matrices $y_{i}$ with entries in $\mathfrak{r}_{v}$. Now the right-hand side belongs to $G L_{N}\left(\mathfrak{r}_{v}\right)$. Therefore, taking the determinant, we obtain (24.8).

If $r=n$, we can take $\beta=1$ and $\zeta=1$. In this case formula (24.7) takes a simpler form:

$$
\begin{equation*}
\left.\mathrm{il}_{\mathfrak{b}}\left(\alpha(q \times \varepsilon h \varepsilon) \sigma^{-1}\right)\right)=\theta^{-n} \operatorname{det}(h) \nu_{\mathfrak{b}}\left(h^{-1} \xi g\right)^{-1} . \tag{24.10}
\end{equation*}
$$

24.4. Lemma. Let $w, w^{\prime} \in \mathcal{H}^{r}$ and $z \in \mathcal{H}^{n}$; let $\xi \in G^{r}$ and $\beta \in G^{n}$; further let $\tau_{\nu}$ be defined by (24.1). Then we have the following formulas:

$$
\begin{align*}
& j\left(\tau_{\nu}, \operatorname{diag}[z, w]\right)=\operatorname{det}\left[1_{\nu}-\wp_{\nu}(w) \wp_{\nu}(z)\right]  \tag{24.11a}\\
& j\left(\tau_{r}, \operatorname{diag}[z, w]\right)=\operatorname{det}\left[1_{r}-w \cdot \wp_{r}(z)\right]=j\left(\eta_{r}^{-1}, w\right) \operatorname{det}\left[\eta_{r}^{-1} w+\wp_{r}(z)\right] \\
& j\left(\tau_{r}\left(\left(\xi \times 1_{2 t}\right) \beta \times 1_{2 r}\right), \operatorname{diag}[z, w]\right) \\
& \quad=j_{\beta}(z) j_{\xi}\left(\wp_{r}(\beta z)\right) j\left(\eta_{r}^{-1}, w\right) \operatorname{det}\left[\eta_{r}^{-1} w+\xi \wp_{r}(\beta z)\right]
\end{align*}
$$

Proof. The first two formulas can be verified by a straightforward calculation. Now, ptting $y=\wp_{r}(\beta z)$, we have $\wp_{r}((\xi \times 1) \beta z)=\xi y$ and

$$
\begin{aligned}
& j\left(\tau_{r}\left(\left(\xi \times 1_{2 t}\right) \beta \times 1_{2 r}\right), \operatorname{diag}[z, w]\right) \\
& \quad=j\left(\tau_{r}, \operatorname{diag}[(\xi \times 1) \beta z, w]\right) j\left(\left(\xi \times 1_{2 t}\right) \beta \times 1, \operatorname{diag}[z, w]\right) \\
& \quad=\operatorname{det}[1-w \cdot(\xi y)] j_{\xi \times 1}(\beta z) j_{\beta}(z)=\operatorname{det}\left[\eta^{-1} w+\xi y\right] j\left(\eta^{-1}, w\right) j_{\xi}(y) j_{\beta}(z)
\end{aligned}
$$

by (23.4) and (24.11b). This proves (24.12).
We now define functions $\delta(w)$ and $\delta\left(w^{\prime}, w\right)$ for $w, w^{\prime} \in \mathcal{H}^{r}$ by

$$
\begin{equation*}
\delta(w)=\delta(w, w), \quad \delta\left(w^{\prime}, w\right)=\left(\operatorname{det}\left[(i / 2)\left(w^{*}-w^{\prime}\right)_{v}\right]\right)_{v \in \mathbf{a}} \quad\left(\in \mathbf{C}^{\mathbf{a}}\right) \tag{24.13}
\end{equation*}
$$

The symbol $\delta(w)$ is consistent with $\delta$ of (16.36). By [S97, (6.6.9) and (7.14.7)],

$$
\begin{equation*}
\delta\left(\gamma w^{\prime}, \gamma w\right)^{m} \overline{j_{\gamma}^{k}(w)} j_{\gamma}^{k}\left(w^{\prime}\right)=\delta\left(w^{\prime}, w\right)^{m} \text { for every } \gamma \in G^{r} \tag{24.14}
\end{equation*}
$$

if $k$ and $m$ are as in $\S 23.2$, since $j_{\gamma}^{k}=j_{\gamma}^{m} \prod_{v \in \mathbf{a}} \operatorname{det}(\alpha)_{v}^{-k_{v \rho}}$.
24.5. Lemma. Let $h \in \mathbf{R}^{\mathbf{a}}$ and $\mathbf{s}=\left(s_{v}\right)_{v \in \mathbf{a}} \in \mathbf{C}^{\mathbf{a}}$; suppose that $h_{v} \geq 0, \operatorname{Re}\left(s_{v}\right)$ $\geq 0$, and $\operatorname{Re}\left(s_{v}\right)+\left(h_{v} / 2\right)>2 \lambda_{r}-1$ for every $v \in \mathbf{a}$, where $\lambda_{r}=(r+1) / 2$ in Case $S P$ and $\lambda_{r}=r$ in Case UT. Then for every holomorphic function $f$ on $\mathcal{H}^{r}$ such that $\delta(w)^{h / 2} f(w)$ is bounded, we have

$$
\begin{gathered}
c_{h}(\mathbf{s}) \delta\left(w^{\prime}\right)^{-\mathbf{s}} f\left(w^{\prime}\right)=\int_{\mathcal{H}^{r}} \delta\left(w^{\prime}, w\right)^{-h}\left|\delta\left(w^{\prime}, w\right)\right|^{-2 \mathbf{s}} \delta(w)^{h+\mathbf{s}} f(w) \mathbf{d} w \\
\text { with } \quad c_{h}(\mathbf{s})=2^{a} \pi^{b} \prod_{v \in \mathbf{a}} \Gamma_{r}^{\iota}\left(s_{v}+h_{v}-\lambda_{r}\right) \Gamma_{r}^{\iota}\left(s_{v}+h_{v}\right)^{-1}
\end{gathered}
$$

where $\mathbf{d} w=\prod_{v \in \mathbf{a}} \mathbf{d} w_{v}$ with $\mathbf{d} w_{v}$ defined as in Lemma 3.4, $a=r(r+3)[F: \mathbf{Q}] / 2$ in Case $S P, a=2 r^{2}[F: \mathbf{Q}]$ in Case $U T, b=r \lambda_{r}[F: \mathbf{Q}]$ in both cases, and $\Gamma_{r}^{\iota}$ is defined by (16.47) with $\iota=[K: F]$.

This is a restatement of [S97, Propositions A2.9 and A2.11]. The expression for the exponent $a$ can be given uniformly $a=r\left(1+\lambda_{r}\right)[F: \mathbf{Q}]$ in both cases, if we take the measure on $\mathcal{H}_{r}$ in Case UT to be $\operatorname{det}(y)^{-2 r} d x d y$ described in $\S 5.12$.
24.6. We now consider $E_{\mathbf{A}}$ of (16.27) with the present $G^{N}$ as $G$ there. We assume $k$ to be integral; we shall make comments in the case of half-integral $k$ in $\S 25.6$. Recall that $E_{\mathbf{A}}$ is determined by the set of data $\{k, \mathfrak{b}, \mathfrak{c}, \chi\}$ satisfying ( $16.24 \mathrm{a}, \mathrm{b}, \mathrm{c}$ ). We put $m=k$ in Case SP and $m=\left(k_{v \rho}+k_{v}\right)_{v \in \mathbf{a}}$ in Case UT as we did in the previous sections. We are interested in $E_{\mathbf{A}}\left(x q_{1} \sigma^{-1}, s\right)$ with the elements $\sigma$ and $q_{1}$ of $G_{\mathrm{h}}^{N}$ given by (24.5) and $q_{1}=q \times \varepsilon h \varepsilon$ as in Lemma 24.3. To be explicit, we have

$$
\begin{equation*}
E_{\mathbf{A}}\left(x q_{1} \sigma^{-1}, s\right)=\sum_{\alpha \in A} \mu\left(\alpha x q_{1} \sigma^{-1}\right) \varepsilon\left(\alpha x q_{1} \sigma^{-1}\right)^{-s}, \quad A=P^{N} \backslash G^{N} \tag{24.15}
\end{equation*}
$$

The function $\mu$ is given by (16.26a, b) with $P^{N}$ and $D^{N}$ as $P$ and $\widetilde{D}$ there. For some practical reasons, we hereafter use the letter $a$ instead of $h$; thus

$$
\begin{equation*}
q_{1}=q \times \varepsilon a \varepsilon, \quad a \in G_{\mathbf{h}}^{r}, \quad a_{v}=1 \text { for every } v \mid c \tag{24.15a}
\end{equation*}
$$

Now, from $E_{\mathbf{A}}\left(x q_{1} \sigma^{-1}, s\right)$ we obtain a function $H_{q, a}(\mathfrak{z}, s)$ of $(\mathfrak{z}, s) \in \mathcal{H}^{N} \times \mathbf{C}$ by the standard principle, that is, $H_{q, a}(y(\mathbf{i}), s)=E_{\mathbf{A}}\left(q_{1} \sigma^{-1} y, s\right) j_{y}^{k}(\mathbf{i})$ for every $y \in G_{\mathbf{a}}^{N}$. Suppressing the variable $s$, we write the function simply $H_{q, a}(\mathfrak{z})$. Then putting $\mathfrak{z}=y(\mathbf{i})$, we have $H_{q, a}(\mathfrak{z})=\sum_{\alpha \in A} p_{\alpha}(\mathfrak{z})$ with

$$
p_{\alpha}(\mathfrak{z})=\mu\left(\alpha q_{1} \sigma^{-1} y\right) \varepsilon\left(\alpha q_{1} \sigma^{-1} y\right)^{-s} j_{y}^{k}(\mathbf{i})
$$

From (16.23a) and (16.26b) we easily see that

$$
\begin{equation*}
p_{\alpha}(\mathfrak{z})=\mu\left(\alpha_{\mathrm{h}} q_{1} \sigma^{-1}\right) \varepsilon\left(\alpha_{\mathrm{h}} q_{1} \sigma^{-1}\right)^{-s} \delta(\mathfrak{z})^{s \mathbf{a}-(m-i \kappa) / 2} \|_{k} \alpha \tag{24.16}
\end{equation*}
$$

where $\delta(\mathfrak{z})=\left(\operatorname{det}\left((i / 2)\left(\mathfrak{z}^{*}-\mathfrak{z}\right)_{v}\right)\right)_{v \in \mathbf{a}}$. By Lemma 24.1, $P^{N} \backslash G^{N}$ can be given as $\bigsqcup_{\nu=0}^{r} A_{\nu}$ with $A_{\nu}=P^{N} \backslash P^{N} \tau_{\nu}\left(G^{n} \times G^{r}\right)$. Put

$$
\begin{equation*}
\mathcal{E}_{\nu}(\mathfrak{z})=\sum_{\alpha \in A_{\nu}} p_{\alpha}(\mathfrak{z}), \quad \mathcal{E}_{\nu}(z, w)=\mathcal{E}_{\nu}(\operatorname{diag}[z, w]) \quad\left(z \in \mathcal{H}^{n}, w \in \mathcal{H}^{r}\right) \tag{24.17}
\end{equation*}
$$

We have then $H_{q, a}=\sum_{\nu=0}^{r} \mathcal{E}_{\nu}$. The functions $H_{q, a}$ and $\mathcal{E}_{\nu}$ involve $s$, but for the moment we suppress it. From (24.16) we easily see that

$$
\begin{equation*}
p_{\alpha} \| \alpha^{\prime}=p_{\alpha \alpha^{\prime}} \text { for every } \alpha^{\prime} \in G^{N} \cap D^{\prime} \tag{24.18}
\end{equation*}
$$

with a suitable open subgroup $D^{\prime}$ of $D^{N}$ independent of $\alpha$. Take a congruence subgroup $\Gamma$ of $G^{r}$ such that $1_{2 n} \times \Gamma \subset D^{\prime}$. Then, from (24.17) and (24.18) we obtain

$$
\begin{equation*}
\mathcal{E}_{\nu} \|\left(1_{2 n} \times \gamma\right)=\mathcal{E}_{\nu} \text { for every } \gamma \in \Gamma . \tag{24.19}
\end{equation*}
$$

Our next task is to obtain an explicit form of $\mathcal{E}_{r}(z, w)$. Given $y \in G_{\mathbf{a}}^{N}$ and $\alpha \in A_{r}$, we have $\mu\left(\alpha y q_{1} \sigma^{-1}\right) \neq 0$ if and only if $\alpha y q_{1} \sigma^{-1} \in P_{\mathbf{A}}^{N} D^{N}$. Since $\left(q_{1}\right)_{v}=$ 1 for every $v \mid \mathfrak{c}$ and $P_{v}^{N} D_{v}^{N}=G_{v}$ for $v \nmid \mathfrak{c}$, we have $P_{\mathbf{A}}^{N} D^{N} \sigma q_{1}^{-1}=P_{\mathbf{A}}^{N} D^{N} \sigma$. Therefore by Lemma 24.2 (2) we can replace $A_{r}$ by the set of elements $\alpha$ of the form $\alpha=\tau_{r}((\xi \times 1) \beta \times 1)$ with $\xi \in \Xi$ and $\beta \in \mathcal{R}_{q}$; we have $\beta \in \zeta q C^{n, r} q^{-1}$ where $\zeta$ and $q$ are as in (23.25) and (23.27). Changing the notation, we use the letter $b$ instead of $g$ in (23.25); thus $q=b \times \operatorname{diag}[\widehat{\varphi}, \varphi]$ with $b \in G_{\mathrm{h}}^{r}, b_{v}=1$ for every $v \mid$ c. Put $x=\alpha y q_{1} \sigma^{-1}$ with such an $\alpha$ and $y \in G_{\mathbf{a}}^{N}$ as above. Since $x \in P_{\mathbf{A}}^{N} D^{N}$, we can put $x=p w$ with $p \in P_{\mathbf{A}}^{N}$ and $w \in D^{N}$. Then we have $d_{x}=d_{p} d_{w}$ and $\operatorname{det}\left(d_{p}\right) \mathfrak{r}=\mathrm{il}_{\mathfrak{b}}(x)$, and so by (24.7) and (24.8),

$$
\begin{aligned}
& \chi_{\mathbf{h}}\left(\operatorname{det}\left(d_{p}\right)\right)^{-1} \chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{w}\right)\right)^{-1}=\left(\chi_{\mathbf{h}} / \chi_{\mathfrak{c}}\right)\left(\operatorname{det}\left(d_{p}\right)\right)^{-1} \chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{x}\right)\right)^{-1} \\
& \quad=\chi^{*}\left(\nu_{\mathfrak{b}}\left(a^{-1} \xi e b\right)\right)\left(\chi_{\mathbf{h}} / \chi_{\mathfrak{c}}\right)\left(\theta^{-r} \operatorname{det}(a) \operatorname{det}(s \varphi)\right)^{-1} \chi_{\mathfrak{c}}\left(\theta^{r} \lambda_{r}(\beta)^{-1}\right) \\
& \quad=\chi_{\mathbf{h}}\left(\theta^{-r} \operatorname{det}(a) \operatorname{det}(\varphi)\right)^{-1} \chi^{*}\left(\nu_{\mathfrak{b}}\left(a^{-1} \xi e b\right)\right) \chi_{\mathbf{a}}\left(\lambda_{r}(\beta)\right)\left(\chi_{\mathbf{h}} / \chi_{\mathfrak{c}}\right)\left(\operatorname{det}(s)^{-1} \lambda_{r}(\beta)\right) .
\end{aligned}
$$

24.7. To simplify our notation, we put $\eta=\eta_{r}, \wp=\wp_{r}$, and

$$
\begin{align*}
& \mathbf{s}=s \mathbf{a}-(m-i \kappa) / 2 \quad\left(\in \mathbf{C}^{\mathbf{a}}\right)  \tag{24.20a}\\
& \nu_{1}(\xi)=\nu_{6}\left(a^{-1} \xi e b\right), \quad N_{1}(\xi)=N\left(\nu_{1}(\xi)\right) \quad(\xi \in \Xi)
\end{align*}
$$

where $N(\mathfrak{r})$ denotes the absolute norm of an $\mathfrak{r}$-ideal $\mathfrak{x}$. These are temporary and will be discarded soon. It should be remembered that $\nu_{1}$ depends on $\mathfrak{b}, a, b$, and $e$, among which $a, b$, and $\mathfrak{b}$ can be fixed, but $e$ depends on $\beta$. Now, by (16.23a) and (24.7),

$$
\varepsilon\left(\alpha_{\mathbf{h}} q_{1} \sigma^{-1}\right)=\varepsilon\left(x_{\mathbf{h}}\right)=N\left(\mathrm{il}_{\mathfrak{b}}(x)\right)^{-u}=N_{1}(\xi)^{u}\left|\theta^{-r} \operatorname{det}(s \varphi)\right|_{K}^{u},
$$

where $u=2 /[K: F]$. (Notice that $|\operatorname{det}(a)|_{K}=1$, since $\operatorname{det}(a) \operatorname{det}(a)^{\rho}=1$.) Since $\operatorname{det}(s) \mathfrak{r}=\mathfrak{a}_{r}^{q}(\beta)$, the last factor of the equality at the end of $\S 24.6$ can be written $\chi^{*}\left(\lambda_{r}(\beta) \mathfrak{a}_{r}^{q}(\beta)^{-1}\right)$. Combining all these formulas with (16.26b) and (24.16), we obtain

$$
\begin{align*}
p_{\alpha}(\mathfrak{z})=\chi_{\mathbf{h}} & \left(\theta^{-r} \operatorname{det}(a) \operatorname{det}(\varphi)\right)^{-1}\left|\theta^{-r} \operatorname{det}(\varphi)\right|_{K}^{-u s}  \tag{24.21}\\
& \cdot \chi^{*}\left(\nu_{1}(\xi)\right) N_{1}(\xi)^{-u s} N\left(\mathfrak{a}_{r}^{q}(\beta)\right)^{u s} \chi[\beta]\left(\delta^{\mathbf{s}} \|_{k} \alpha\right)(\mathfrak{z}),
\end{align*}
$$

where $\chi[\beta]=\chi_{\mathbf{a}}\left(\lambda_{r}(\beta)\right) \chi^{*}\left(\lambda_{r}(\beta) \mathfrak{a}_{r}^{q}(\beta)^{-1}\right)$. Put

$$
\begin{align*}
& M\left(w^{\prime}, w\right)=\delta(w)^{m+\mathbf{s}} \operatorname{det}\left(w^{\prime}-w^{*}\right)^{-m}\left|\operatorname{det}\left(w^{\prime}-w^{*}\right)\right|^{-2 \mathbf{s}} \quad\left(w, w^{\prime} \in \mathcal{H}^{r}\right)  \tag{24.22a}\\
& \mathcal{A}(s)=\chi\left(\theta^{-r} \operatorname{det}(a) \operatorname{det}(\varphi)\right)^{-1}\left|\theta^{-r} \operatorname{det}(\varphi)\right|_{K}^{-u s} . \tag{24.22b}
\end{align*}
$$

Now $\left(\delta^{\mathbf{s}} \|_{k} \alpha\right)(\mathfrak{z})=j_{\alpha}^{k}(\mathfrak{z})^{-1}\left|j_{\alpha}(\mathfrak{z})\right|^{-2 \mathbf{s}} \delta(\mathfrak{z})^{\mathbf{s}}$. By (24.12) and (24.13) we obtain

$$
\begin{aligned}
& \left(\delta^{\mathbf{s}} \|_{k} \alpha\right)(\operatorname{diag}[z, w])=j_{\xi}^{k}(\wp(\beta z))^{-1}\left|j_{\xi}(\wp(\beta z))\right|^{-2 \mathbf{s}}\left(\delta(w)^{\mathbf{s}} \|_{k} \eta^{-1}\right) \\
& \quad \cdot \operatorname{det}\left[\xi \wp(\beta z)+\eta^{-1} w\right]^{-m}\left|\operatorname{det}\left[\xi \wp(\beta z)+\eta^{-1} w\right]\right|^{-2 \mathbf{s}}\left(\delta(z)^{\mathbf{s}} \|_{k} \beta\right)
\end{aligned}
$$

Transforming our functions by $\eta$, we put

$$
\begin{equation*}
\mathcal{F}(z, w)=H_{q, a}(\operatorname{diag}[z, \eta w]) j_{\eta}^{k}(w)^{-1}, \quad \mathcal{F}_{\nu}(z, w)=\mathcal{E}_{\nu}(z, \eta w) j_{\eta}^{k}(w)^{-1} \tag{24.23}
\end{equation*}
$$

Change $w$ for $-w^{*}$; then from the above calculations we obtain

$$
\begin{align*}
& \delta(w)^{m} \mathcal{F}_{r}\left(z,-w^{*}\right)=\mathcal{A}(s) \sum_{\beta \in \mathcal{R}_{q}} \sum_{\xi \in \Xi} \chi^{*}\left(\nu_{1}(\xi)\right) N_{1}(\xi)^{-u s} j_{\xi}^{k}(\wp(\beta z))^{-1}  \tag{24.24}\\
& \cdot\left|j_{\xi}(\wp(\beta z))\right|^{-2 \mathbf{s}} \chi[\beta] N\left(\mathfrak{a}_{r}^{q}(\beta)\right)^{u s} M(\xi \wp(\beta z), w)\left(\delta(z)^{\mathbf{s}} \| \beta\right) .
\end{align*}
$$

We note here an easy fact: If $f \in \mathcal{S}_{k}^{r}(\Gamma)$ and $\gamma \in \Gamma$, then the expression

$$
M\left(\xi w^{\prime}, w\right) j_{\xi}^{k}\left(w^{\prime}\right)^{-1}\left|j_{\xi}\left(w^{\prime}\right)\right|^{-2 \mathbf{s}} f(w)
$$

is invariant under $(w, \xi) \mapsto(\gamma w, \gamma \xi)$ for every $\gamma \in \Gamma$. This follows immediately from (24.14).
24.8. We take $\mathcal{B}$ as in (20.5) with $r$ in place of $n$, which we denoted by $\mathcal{B}^{r}$ in $\S 23.4$; we assume that $c_{v}=1$ for every $c \in \mathcal{B}$ and every $v \mid c ;$ we take $\mathcal{B}=\{1\}$ in Case SP; also, in Case UT we take each $c$ in the form $c=\operatorname{diag}[\widehat{d}, d]$ with $d \in G L_{r}(K)_{\mathbf{h}}$ (see [S97, Lemma 9.8 (3)].) We put $\Gamma^{c}=G^{r} \cap c C^{r} c^{-1}$ for each $c \in \mathcal{B}$. Let $\mathbf{f} \in \mathcal{S}_{k}^{r}\left(C^{r}\right)$. For each $c \in \mathcal{B}$ define $f_{c}$ as in (20.3b). For every function $\psi$ on $\mathcal{H}^{r}$ we define a function $\psi^{\sim}=\widetilde{\psi}$ on $\mathcal{H}^{r}$ by $\widetilde{\psi}(w)=\overline{\psi\left(-w^{*}\right)}$. Notice that $\widetilde{\psi} \in \mathcal{S}_{k \rho}^{r}\left(\Gamma^{c}\right)$ if $\psi \in \mathcal{S}_{k}^{r}\left(\Gamma^{c}\right)$. This is easy in Case SP; we need (5.34) in Case UT.

We now assume that $n>r$, and consider

$$
\begin{equation*}
\int_{\mathfrak{D}} \mathcal{F}(z, w) \overline{\psi_{a}(w)} \delta(w)^{m} \mathbf{d} w \quad\left(\mathfrak{D}=\Gamma_{0} \backslash \mathcal{H}^{r}\right) \tag{24.25}
\end{equation*}
$$

with a congruence subgroup $\Gamma_{0}$ of $G^{r}$ contained in $\Gamma^{a}$, where $\psi_{a}=\left(f_{a}\right)^{\sim}$ and $\mathbf{d} w$ is as in Lemma 24.5. This can be written $\sum_{\nu=0}^{r} I_{\nu}$ with

$$
\begin{equation*}
I_{\nu}=\int_{\mathcal{D}} \mathcal{F}_{\nu}(z, w) \overline{\psi_{a}(w)} \delta(w)^{m} \mathbf{d} w \tag{24.26}
\end{equation*}
$$

By (24.19), these integrals are (at least formally) meaningful for a suitable choice of $\Gamma_{0}$. As explained in [S97, §22.12], integrals of this type converge locally uniformly on $\mathcal{H}^{n} \times\left\{s \in \mathbf{C} \mid \operatorname{Re}(s)>\sigma_{0}\right\}$ for some $\sigma_{0} \in \mathbf{R}$. We shall also show in $\S 24.10$ that $I_{\nu}=0$ for $0 \leq \nu<r$.

To compute $I_{r}$, we first note

$$
I_{r}=\int_{\mathfrak{D}} \mathcal{F}_{r}\left(z,-w^{*}\right) f_{a}(w) \delta(w)^{m} \mathbf{d} w
$$

Now $\nu_{1}(\xi)$ and $N_{1}(\xi)$ depend only on $\Gamma^{a} \xi$, and so, in view of the remark at the end of $\S 24.7$, each term of (24.24) times $f_{a}(w)$ is invariant under $(w, \xi) \mapsto(\gamma w, \gamma \xi)$ for every $\gamma \in \Gamma^{a}$. Let $R^{a}$ be a complete set of representatives for $\Gamma^{a} \backslash \Xi$, and let $\mathcal{F}^{\prime}\left(z,-w^{*}\right)$ denote the function such that $\delta(w)^{m} \mathcal{F}^{\prime}\left(z,-w^{*}\right)$ is the right-hand side of (24.24) with $\sum_{\xi \in R^{a}}$ in place of $\sum_{\xi \in \Xi}$. Then

$$
\mathcal{F}_{r}\left(z,-w^{*}\right) \delta(w)^{m} f_{a}(w)=\sum_{\gamma \in \Gamma^{a}} \mathcal{F}^{\prime}\left(z,-(\gamma w)^{*}\right)\left(\delta^{m} f_{a}\right)(\gamma w)
$$

so that

$$
I_{r}=\mu \int_{\mathcal{H}^{r}} \mathcal{F}^{\prime}\left(z,-w^{*}\right) \delta(w)^{m} f_{a}(w) \mathbf{d} w
$$

where $\mu=\left[\Gamma^{a}: \Gamma_{0}\right]\left[\Gamma_{0} \cap \mathfrak{r}^{\times}: 1\right]$. For $c \in G_{\mathrm{h}}^{r}$ put

$$
\begin{equation*}
f_{a} \mid \mathfrak{T}_{c}^{a}=\sum_{\xi \in R^{a}} \chi^{*}\left(\nu_{\mathfrak{b}}\left(a^{-1} \xi c\right)\right) N\left(\nu_{\mathfrak{b}}\left(a^{-1} \xi c\right)\right)^{-u s} f_{a} \|_{k} \xi \tag{24.27}
\end{equation*}
$$

Then $f_{a} \mid \mathfrak{T}_{e b}^{a}=\sum_{\xi \in R^{a}} \chi^{*}\left(\nu_{1}(\xi)\right) N_{1}(\xi)^{-u s} f_{a} \|_{k} \xi$. By Lemma 24.5 , termwise integration yields

$$
\begin{aligned}
I_{r}= & \mu c_{m}(\mathbf{s}) \mathcal{A}(s) \sum_{\beta} \sum_{\xi \in R^{a}} \chi^{*}\left(\nu_{1}(\xi)\right) N_{1}(\xi)^{-u s} \\
& \cdot\left(f_{a} \| \xi\right)(\wp(\beta z)) \delta(\wp(\beta z))^{-\mathbf{s}} \chi[\beta] N\left(\mathfrak{a}_{r}^{q}(\beta)\right)^{u s} \delta^{\mathbf{s}} \| \beta \\
= & \mu c_{m}(\mathbf{s}) \mathcal{A}(s) \sum_{\beta}\left(f_{a} \mid \mathfrak{T}_{e b}^{a}\right)(\wp(\beta z)) \chi[\beta] N\left(\mathfrak{a}_{r}^{q}(\beta)\right)^{u s}(\delta /(\delta \circ \wp))^{\mathbf{s}} \| \beta \\
= & \mu c_{m}(\mathbf{s}) \mathcal{A}(s) \sum_{\beta} \chi[\beta] N\left(\mathfrak{a}_{r}^{q}(\beta)\right)^{u s} \delta\left(z, s ; f_{a} \mid \mathfrak{T}_{e b}^{a}, \kappa\right) \| \beta
\end{aligned}
$$

with $c_{m}$ of Lemma 24.5 and $\delta(\cdots)$ of (23.28). Put

$$
\begin{equation*}
\left.\mathbf{f}\right|_{\chi} \mathfrak{T}=\sum_{\tau \in C \backslash \mathfrak{X} / C}(\mathbf{f} \mid C \tau C) \chi^{*}\left(\nu_{\mathfrak{b}}(\tau)\right) N\left(\nu_{\mathfrak{b}}(\tau)\right)^{-u s} \quad\left(C=C^{r}\right) \tag{24.28}
\end{equation*}
$$

This is obtained from (20.11) by substituting $\chi^{*}\left(\nu_{\mathfrak{b}}(\tau)\right) N\left(\nu_{\mathfrak{b}}(\tau)\right)^{-u s}$ for $\left[\nu_{\mathfrak{b}}(\tau)\right]$. Let $\left.\mathbf{f}\right|_{\chi} \mathfrak{T} \leftrightarrow\left(f_{b}^{\prime}\right)_{b \in \mathcal{B}}$. By Lemma 20.10 we have $f_{e b}^{\prime}=\sum_{a \in \mathcal{B}} f_{a} \mid \mathfrak{T}_{e b}^{a}$.

Define $E\left(x, s ;\left.\mathbf{f}\right|_{\chi} \mathfrak{T}, \chi, C^{n, r}\right)$ by (23.24) with $\left.\mathbf{f}\right|_{\chi} \mathfrak{T}$ in place of $\mathbf{f} ;$ associate a function $E_{q}$ to this by (23.26) for $q$ as in (23.25), and denote it by $E_{q}\left(z, s ;\left.\mathbf{f}\right|_{\chi} \mathfrak{T}\right)$. Combining the above calculation with Proposition 23.5, we obtain a fundamental formula

$$
\begin{align*}
& \mu c_{m}(\mathbf{s}) \chi(\theta)^{r} N(\mathfrak{b r})^{-r u s} E_{q}\left(z, s ; \mathbf{f}_{\chi} \mathfrak{T}\right)  \tag{24.29}\\
& \quad=\sum_{a \in \mathcal{B}} \chi_{\mathbf{h}}(\operatorname{det}(a)) \int_{\mathfrak{D}} J_{q, a}\left(z,-w^{*} ; s\right) f_{a}(w) \delta(w)^{m} \mathbf{d} w
\end{align*}
$$

where $J_{q, a}(z, w ; s)=H_{q, a}(\operatorname{diag}[z, \eta w]) j_{\eta}^{k}(w)^{-1}$. (We can take $\mu$ - to be the same for all $a \in \mathcal{B}$, since $\left[\Gamma^{a}: \Gamma_{0}\right.$ ] does not depend on $a$; see [S97, Lemma 8.15].) Now, our termwise integration can be justified for sufficiently large $\operatorname{Re}(s)$, because of the validity of Lemma 24.5 for such an $s$ and of the convergence of (24.25) and the series expressing $E_{q}$ in (23.26). Assuming that $\mathbf{f}$ is an eigenform, multiply both sides of (24.29) by $\Lambda_{\mathrm{c}}^{n+r}(s, \chi)$. Then we obtain

$$
\begin{align*}
& \mu c_{m}(\mathbf{s}) \chi(\theta)^{r} N(\mathfrak{b r})^{-r u s} \mathcal{F}_{q}(z, s ; \mathbf{f}, \chi, C)  \tag{24.29a}\\
& =\sum_{a \in \mathcal{B}} \chi_{\mathbf{h}}(\operatorname{det}(a)) \int_{\mathfrak{D}} \Lambda_{\mathfrak{c}}^{n+r}(s, \chi) J_{q, a}\left(z,-w^{*} ; s\right) f_{a}(w) \delta(w)^{m} \mathbf{d} w .
\end{align*}
$$

24.9. So far we have assumed that $n>r$ in the above treatment, but everything is meaningful even if $n=r$. Suppose $n=r$; then the sum $\sum_{\beta}$ in $\S 24.8$ consists of a single term for $\beta=1 ; \wp(z)=z$ for $z \in \mathcal{H}^{n}=\mathcal{H}^{r}, q=b, \zeta=1, e=1$, and so we have $I_{r}=\mu c_{m}(\mathbf{s}) \mathcal{A}(s) f_{a} \mid \mathfrak{T}_{b}^{a}$. Then (24.29) can be written

$$
\begin{align*}
& \mu c_{m}(\mathbf{s}) \chi(\theta)^{n} N(\mathbf{b r})^{-n u s} f_{b}^{\prime}(z, s)  \tag{24.30}\\
& \quad=\sum_{a \in \mathcal{B}} \chi_{\mathbf{h}}(\operatorname{det}(a)) \int_{\mathcal{D}} J_{b, a}\left(z,-w^{*} ; s\right) f_{a}(w) \delta(w)^{m} \mathbf{d} w
\end{align*}
$$

where $f_{b}^{\prime}(z, s)$ is the $b$-component of $\mathbf{f} \mid \chi \mathfrak{T}$. In particular, if $\mathbf{f}$ is an eigenfunction as in $\S 20.6$, we have $f_{b}^{\prime}(z, s)=\mathfrak{T}(u s, \mathbf{f}, \chi) f_{b}(z)$.
24.10. Let us now prove that $I_{\nu}=0$ for $0 \leq \nu<r$. Fixing $\nu$, take a complete set of representatives $\mathcal{R}^{n}$ (resp. $\mathcal{R}^{r}$ ) for $P^{n, \nu} \backslash G^{n}$ (resp. $P^{r, \nu} \backslash G^{r}$ ), and put $p_{\xi, \beta, \gamma}=p_{\alpha}$ for $\alpha=\tau_{\nu}\left(\left(\xi \times 1_{2 n-2 \nu}\right) \beta \times \gamma\right)$ with $\xi \in G^{\nu}, \beta \in G^{n}$, and $\gamma \in G^{r}$. By Lemma 24.1, $\mathcal{E}_{\nu}=\sum_{\xi \in G^{\nu}} \sum_{\beta \in \mathcal{R}^{n}} \sum_{\gamma \in \mathcal{R}^{r}} p_{\xi, \beta, \gamma}$. Let $\Gamma$ be a congruence subgroup of $G^{r}$ such that $1_{2 n} \times \Gamma \subset D^{\prime}$ with $D^{\prime}$ as in (24.18). Take a finite subset $T$ of $G^{r}$ so that $G^{r}=\bigsqcup_{\tau \in T} P^{r, \nu} \tau \Gamma$ (see [S97, Lemma 9.8]). Then we can take $\mathcal{R}^{r}=\bigsqcup_{\tau \in T} \mathcal{R}_{\tau}$ with $\mathcal{R}_{\tau}=\left(P^{r, \nu} \cap \tau \Gamma \tau^{-1}\right) \backslash \tau \Gamma$. Put $g_{\tau}=\sum_{\xi \in G^{\nu}} \sum_{\beta \in \mathcal{R}^{n}} p_{\xi, \beta, \tau} \|\left(1 \times \tau^{-1}\right)$. If $\varepsilon \in \tau \Gamma$, then $1 \times \tau^{-1} \varepsilon \in D^{\prime}$, and so $p_{\xi, \beta, \tau} \|\left(1 \times \tau^{-1} \varepsilon\right)=p_{\xi, \beta, \varepsilon}$ by (24.18). Thus $\mathcal{E}_{\nu}=\sum_{\tau \in T} \sum_{\varepsilon \in \mathcal{R}_{\tau}} g_{\tau} \|(1 \times \varepsilon)$. Now, given $\gamma \in P^{r, \nu} \cap \tau \Gamma \tau^{-1}$, put $\omega=\kappa_{\nu} \pi_{\nu}(\gamma)^{-1} \kappa_{\nu}$ with $\kappa_{\nu}$ of Lemma 24.1. Then by that lemma we have

$$
\tau_{\nu}\left(\left(\xi \times 1_{2 n-2 \nu}\right) \beta \times \tau\right)\left(1_{2 n} \times \tau^{-1} \gamma \tau\right) \in P_{N} \tau_{\nu}\left(\left(\omega \xi \times 1_{2 n-2 \nu}\right) \beta \times \tau\right)
$$

and hence $p_{\xi, \beta, \tau} \|\left(1 \times \tau^{-1} \gamma \tau\right)=p_{\omega \xi, \beta, \tau}$, since $p_{\alpha}$ depends only on $P^{N} \alpha$. This shows that $g_{\tau} \|(1 \times \gamma)=g_{\tau}$ for such a $\gamma$. Next, put

$$
\begin{equation*}
p^{\circ}(z, w)=p(\operatorname{diag}[z, w]) \quad\left(z \in \mathcal{H}^{n}, w \in \mathcal{H}^{r}\right) \tag{24.31}
\end{equation*}
$$

for a function $p$ on $\mathcal{H}^{N}$; let $R$ be the unipotent radical of $P^{r, \nu}$ and let $\zeta \in R_{\mathbf{a}}$. Since $\wp_{\nu}(\zeta w)=\wp_{\nu}(w)$ for every $w \in \mathcal{H}^{r}$ and $j_{\zeta}=1$, from (24.11a) we see that $\left[\delta^{\mathbf{s}} \|\left(\tau_{\nu}\left(1_{2 n} \times \zeta\right)(\beta \times \gamma)\right)\right]^{\circ}=\left[\delta^{\mathbf{s}} \| \tau_{\nu}(\beta \times \gamma)\right]^{\circ}$. Now, for $\alpha=\tau_{\nu} \alpha^{\prime}$ with $\alpha^{\prime}=$ $\left(\xi \times 1_{2 n-2 \nu}\right) \beta \times \tau$ we have $p_{\alpha}\left\|\left(1 \times \tau^{-1}\right)=c_{\alpha} \delta^{s}\right\|\left(\tau_{\nu} \alpha^{\prime}\left(1 \times \tau^{-1}\right)\right)$ with a constant $c_{\alpha}$, and hence $\left[p_{\alpha} \|\left(1 \times \tau^{-1}\right)(1 \times \zeta)\right]^{\circ}=c_{\alpha}\left[\delta^{s} \|\left(\tau_{\nu}(1 \times \zeta) \alpha^{\prime}\left(1 \times \tau^{-1}\right)\right)\right]^{\circ}=c_{\alpha}\left[\delta^{\mathbf{s}} \|\left(\tau_{\nu} \alpha^{\prime}(1 \times\right.\right.$ $\left.\left.\left.\tau^{-1}\right)\right)\right]^{\circ}=\left[p_{\alpha} \|\left(1 \times \tau^{-1}\right)\right]^{\circ}$. Thus $\left[g_{\tau} \|(1 \times \zeta)\right]^{\circ}=\left(g_{\tau}\right)^{\circ}$ for every $\zeta \in R_{\mathbf{a}}$. Put $g(w)=\left(g_{\tau}\right)^{\circ}(z, w)$ with fixed $\tau$ and $z$. We have shown that $g \| \gamma=g$ for every $\gamma \in P^{r, \nu} \cap \tau \Gamma \tau^{-1}$ and $g \| \zeta=g$ for every $\zeta \in R_{\mathbf{a}}$. Therefore, by [S97, Lemma A3.8], $\left\langle\psi, \sum_{\varepsilon \in \mathcal{R}_{\tau}} g \| \varepsilon\right\rangle=0$ for every $\psi \in \mathcal{S}_{k}^{r}$ if $r>\nu$. Consequently $\left\langle\psi(w), \mathcal{E}_{\nu}(z, w)\right\rangle=$ 0 for every such $\psi$, at least for sufficiently large $\operatorname{Re}(s)$. Tranforming $\mathcal{E}_{\nu}$ by $\eta$, we obtain $I_{\nu}=0$.
24.11. Lemma. Let $k$ be an integral or a half-integral weight, and $\Psi$ a subfield of $\mathbf{C}$ containing the Galois closure of $K$ over $\mathbf{Q}$. In order to emphasize the dimensionality, denote by $\mathcal{N}_{k}^{n, p}(\Psi)\left(r e s p . \mathcal{N}_{k}^{n, p}(\Gamma, \Psi)\right)$ the set $\mathcal{N}_{k}^{p}(\Psi)\left(r e s p . \mathcal{N}_{k}^{p}(\Gamma, \Psi)\right)$ of $\S 14.11$ defined with respect to $G^{n}$. If $f \in \mathcal{N}_{k}^{n+r, p}(\Psi)$, then $f^{\circ}(z, w)$ can be written as a finite sum $f^{\circ}(z, w)=\sum_{a=1}^{t} g_{a}(z) h_{a}(w)$ with $g_{a} \in \mathcal{N}_{k}^{n, p}(\Psi)$ and $h_{a} \in \mathcal{N}_{k}^{r, p}(\Psi)$. In particular, if $p=0$, the conclusion holds for every subfield $\Psi$ of $\mathbf{C}$ containing the field $\Phi_{k}$ of Theorem 10.4 (5) or Theorem 10.7 (5).

Proof. Take congruence subgroups $\Gamma^{i}$ of $G^{i}$ for $i=n, r$, and $n+r$ so that $\Gamma^{n} \times \Gamma^{r} \subset \Gamma^{n+r}$ and $f \in \mathcal{N}_{k}^{n+r, p}\left(\Gamma^{n+r}, \Psi\right)$. We can take each $\Gamma^{i}$ to be a principal congruence subgroup of some level such that $\mathcal{N}_{k}^{i, p}\left(\Gamma^{i}\right)=\mathcal{N}_{k}^{i, p}\left(\Gamma^{i}, \Psi\right) \otimes_{\Psi} \mathbf{C}$ as in Proposition 14.13 (1). Let $\left\{h_{a}\right\}_{a=1}^{t}$ be a $\Psi$-basis of $\mathcal{N}_{k}^{r, p}\left(\Gamma^{r}, \Psi\right)$. We easily see that $f^{\circ}(z, w)$ as a function of $z$ (resp. $w$ ) belongs to $\mathcal{N}_{k}^{n, p}\left(\Gamma^{n}\right)$ (resp. $\mathcal{N}_{k}^{r, p}\left(\Gamma^{r}\right)$ ). (If $k$ is half-integral, Proposition A2.12 is essential.) Therefore for each fixed $z$ we have $f^{\circ}(z, w)=\sum_{a=1}^{t} g_{a}(z) h_{a}(w)$ with complex numbers $g_{a}(z)$ uniquely determined by $z$ and $a$. Since $\bigcap_{w}\left\{x \in \mathbf{C}^{t} \mid \sum_{a=1}^{t} x_{a} h_{a}(w)\right\}=\{0\}$, we can find $t$ points $w_{1}, \ldots, w_{t}$ of $\mathcal{H}^{r}$ such that $\operatorname{det}\left(h_{a}\left(w_{b}\right)\right)_{a, b=1}^{t} \neq 0$. Solving the linear equations $f^{\circ}\left(z, w_{b}\right)=\sum_{a=1}^{t} g_{a}(z) h_{a}\left(w_{b}\right)$, we find that $g_{a} \in \mathcal{N}_{k}^{n, p}$. Let $\sigma \in \operatorname{Aut}(\mathbf{C} / \Psi)$; then $f^{\sigma}=f$ and $\left(h_{a}\right)^{\sigma}=h_{a}$. Comparing the Fourier coefficients of $f$ with those of $g_{a}$ and $h_{a}$, and applying $\sigma$ to them, we can easily verify that $f^{\circ}(z, w)=$
$\sum_{a=1}^{t}\left(g_{a}\right)^{\sigma}(z) h_{a}(w)$, and hence $\left(g_{a}\right)^{\sigma}=g_{a}$, that is, $g_{a} \in \mathcal{N}_{k}^{n, p}(\Psi)$. This proves our lemma in the general case. In the case $p=0$, we have $\mathcal{M}_{k}\left(\Gamma^{r}\right)=\mathcal{M}_{k}\left(\Gamma^{r}, \Phi_{k}\right) \otimes_{\Phi_{k}}$ C by Theorem 10.4 (5) or Theorem 10.7 (5), and hence the above argument is valid if $\Phi_{k} \subset \Psi$.

## 25. Proof of Theorems in Sections 20 and 23

25.1. Our computation of Section 24 is sufficient for the proof of Theorem 23.9 and a special case of Theorem 20.14. To prove the most general case of the latter theorem, we have to introduce certain differential operators. They are necessary only in the case $n=r$, and so we speak of $G^{2 n}$ instead of $G^{N}$. Returning to the setting of Section 12, we take $T$ and $S_{p}(T)$ as in $\S 13.13$ and also $\mathfrak{K}$ as in (14.4), all with $2 n$ instead of $n$ there. Thus $T_{v}=\mathbf{C}_{2 n}^{2 n}$ in Case UT and $T_{v}=\left\{\left.z \in \mathbf{C}_{2 n}^{2 n}\right|^{t} z=\right.$ $z\}$ in Case $\mathrm{SP} ; \mathfrak{K}=G L_{2 n}(\mathbf{C})^{\mathbf{b}}$. We shall simply write $\tau$ for the representation $\tau^{p}$ of $\S \S 13.13$ and 14.4.

Given a representation $\{\omega, X\}$ of $\mathfrak{K}$, an irreducible subspace $Z$ of $S_{p}(T)$ as in §13.13, and $\zeta \in Z$, we define differential operators $B_{\zeta}$ and $C_{\zeta}$ on $\mathcal{H}^{2 n}$ by

$$
\begin{equation*}
B_{\zeta} f=\left(D_{\omega}^{Z} f\right)(\zeta), \quad C_{\zeta} f=\left(E^{Z} f\right)(\zeta) \tag{25.1}
\end{equation*}
$$

for $f \in C^{\infty}\left(\mathcal{H}^{2 n}, X\right)$, where $D_{\omega}^{Z}$ and $E^{Z}$ are defined in (13.22) and §14.4. Then

$$
\begin{align*}
& \left(B_{\zeta} f\right) \|_{\omega} \alpha=B_{\psi_{1}}\left(f \|_{\omega} \alpha\right) \text { with } \psi_{1}(u)=\zeta\left(\lambda_{\alpha}(\mathfrak{z}) u \cdot{ }^{t} \mu_{\alpha}(\mathfrak{z})\right)  \tag{25.2a}\\
& \left(C_{\zeta} f\right) \|_{\omega} \alpha=C_{\psi_{2}}\left(f \|_{\omega} \alpha\right) \text { with } \psi_{2}(u)=\zeta\left({ }^{t} \lambda_{\alpha}(\mathfrak{z})^{-1} u \mu_{\alpha}(\mathfrak{z})^{-1}\right),
\end{align*}
$$

where $u \in T$ and $\mathfrak{z}$ is the variable on $\mathcal{H}^{2 n}$. Indeed, put $\lambda=\lambda_{\alpha}(\mathfrak{z}), \mu=\mu_{\alpha}(\mathfrak{z}), g=$ $D_{\omega}^{Z} f$, and define $\psi_{1}$ as above for a given $\zeta \in Z$. Then, by (12.21) and (12.24a), or rather by their generalizations mentioned in §13.13, we have

$$
\begin{aligned}
B_{\psi_{1}}\left(f \|_{\omega} \alpha\right) & =\left[D_{\omega}^{Z}\left(f \|_{\omega} \alpha\right)\right]\left(\psi_{1}\right)=\omega(\lambda, \mu)^{-1}(g \circ \alpha)\left(\tau\left({ }^{t} \lambda,{ }^{t} \mu\right)^{-1} \psi_{1}\right) \\
& =\left(g \|_{\omega} \alpha\right)(\zeta)=g(\zeta)\left\|_{\omega} \alpha=\left(B_{\zeta} f\right)\right\|_{\omega} \alpha,
\end{aligned}
$$

which gives (25.2a). The other formula can be proved in the same way.
Let us now consider the case $\omega(x)=\operatorname{det}(x)^{k}$ with $k \in \mathbf{Z}^{\mathbf{b}}$ and $Z=\bigotimes_{v \in \mathbf{a}} Z_{v}$, where $Z_{v}$ is the irreducible subspace of $S_{n e_{v}}\left(T_{v}\right)$ whose highest weight vector is $\operatorname{det}_{n}(u)^{e_{v}}$ for $u \in T_{v}$ with $0 \leq e_{v} \in \mathbf{Z}$ (see Theorem 12.7). We assume $e_{v} \leq 1$ in Case SP. For $u \in \mathbf{C}_{2 n}^{2 n}$ we denote by $u_{\ell}$ the lower left $(n \times n)$-block of $z$. We then define an elements $\varphi$ of $S_{n e}(T)$ by

$$
\begin{equation*}
\varphi(u)=\prod_{v \in \mathbf{a}} \operatorname{det}\left(\left(u_{v}\right)_{\ell}\right)^{e_{v}} \quad(u \in T) \tag{25.3}
\end{equation*}
$$

Then $\varphi \in Z$. This is clear in Case UT, since we can find an element $(a, b) \in \mathfrak{K}$ such that the upper right $(n \times n)$-block of ${ }^{t} a u b$ is $u_{\ell}$. In Case SP , the inclusion $\varphi \in Z$, which holds only under the assumption that $e_{v} \leq 1$, can be seen as follows.

Clearly it is sufficient to consider a single $v$, and so we drop the subscript $v$, Using the symbols $a_{x}, b_{x}, c_{x}$, and $d_{x}$ for $x \in \mathbf{C}_{2 n}^{2 n}$, we have $\left({ }^{t} x u x\right)_{\ell}={ }^{t} b_{x} a_{u} a_{x}+Y$ for $(x, x) \in \mathfrak{K}$ and $u \in T$ with a matrix $Y$ that does not involve $u_{11}$. Thus $\varphi\left({ }^{t} x u x\right)$ as a function of $u$ is of degree $\leq 1$ in $u_{11}$. Let $W$ be the irreducible subspace of $S_{n}\left(T_{v}\right)$ contained in the $\tau(\mathfrak{K})$-span of $\varphi$. A highest weight vector of $W$ can be given in the form $\prod_{\nu} \operatorname{det}_{\nu}(u)^{c_{\nu}}$ with $\left(c_{\nu}\right)$ such that $n=\sum_{\nu} \nu c_{\nu}$. Since this must be of degree $\leq 1$ in $u_{11}$, we see that $W=Z$, which proves that $\varphi \in Z$ in Case SP.
25.2. Given $k \in \mathbf{Z}^{\mathbf{b}}, \beta, \gamma \in G^{n}$, and a function $g$ on $\mathcal{H}^{n} \times \mathcal{H}^{n}$, we define a function $g \|_{k}(\beta \times \gamma)$ on $\mathcal{H}^{n} \times \mathcal{H}^{n}$ by

$$
\begin{equation*}
\left(g \|_{k}(\beta \times \gamma)\right)(z, w)=j_{\beta}^{k}(z)^{-1} j_{\gamma}^{k}(w)^{-1} g(\beta z, \gamma w) \quad\left(z, w \in \mathcal{H}^{n}\right) \tag{25.4}
\end{equation*}
$$

Fixing $k$, let us hereafter write $B_{e}$ and $C_{e}$ for $B_{\varphi}$ and $C_{\varphi}$ with $\varphi$ as above. Employing the symbol $p^{\circ}$ of (24.31) with $r=n$, for $\beta, \gamma \in G^{n}$ we have

$$
\begin{align*}
& {\left[B_{e}\left(f \|_{k}(\beta \times \gamma)\right)\right]^{\circ}=\operatorname{det}(\gamma)^{e}\left(B_{e} f\right)^{\circ} \|_{k+e}(\beta \times \gamma),}  \tag{25.5a}\\
& {\left[C_{e}\left(f \|_{k}(\beta \times \gamma)\right)\right]^{\circ}=\operatorname{det}(\gamma)^{-e}\left(C_{e} f\right)^{\circ} \|_{k-e}(\beta \times \gamma),} \tag{25.5b}
\end{align*}
$$

where we consider $e$ as an element of $\mathbf{Z}^{\mathbf{b}}$ via the natural injection of $\mathbf{Z}^{\mathbf{a}}$ into $\mathbf{Z}^{\mathbf{b}}$. To prove these, take $\omega(x)=\operatorname{det}(x)^{k}$ and $\zeta=\varphi$ in (25.2a). For $\alpha=\beta \times \gamma$ and $\mathfrak{z}=\operatorname{diag}[z, w]$ we have $\left(\lambda_{\alpha}(\mathfrak{z}) u \cdot{ }^{t} \mu_{\alpha}(\mathfrak{z})\right)_{\ell}=\lambda_{\gamma}(w) u_{\ell} \cdot{ }^{t} \mu_{\beta}(z)$, and hence $\psi_{1}=$ $j_{\gamma}^{e \rho}(w) j_{\beta}^{e}(z) \varphi$, which gives (25.5a). Formula (25.5b) follows from (25.2b) in a similar way. $C_{e}$ can be defined in both cases, but actually it is unnecessary in Case SP.

Now take two elements $e$ and $e^{\prime}$ of $\mathbf{Z}^{\mathbf{a}}$ such that $e_{v} \geq 0, e_{v}^{\prime} \geq 0$, and $e_{v} e_{v}^{\prime}=0$ for every $v \in \mathbf{a}$. Put $D_{e, e^{\prime}}=B_{e} C_{e^{\prime}}$. Then from (25.5a, b) we obtain immediately

$$
\begin{equation*}
\left[D_{e, e^{\prime}}\left(f \|_{k}(\beta \times \gamma)\right)\right]^{\circ}=\operatorname{det}(\gamma)^{e-e^{\prime}}\left(D_{e, e^{\prime}} f\right)^{\circ} \|_{k+e-e^{\prime}}(\beta \times \gamma) . \tag{25.6}
\end{equation*}
$$

25.3. Lemma. Let $\tau_{n}$ be as in Lemma 24.1 with $N=2 n$ and let $\mathbf{s}=\left(s_{v}\right)_{v \in \mathbf{a}} \in$ $\mathbf{C}^{\mathbf{a}}$; let $m=k$ in Case SP and $m=\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$ in Case UT. Then

$$
\begin{aligned}
& {\left[D_{e, e^{\prime}}\left(\delta^{\mathbf{s}} \|_{k} \tau_{n}\right)\right]^{\circ}=\Psi(\mathbf{s})\left(\delta^{\mathbf{s}+e^{\prime}} \|_{h} \tau_{n}\right)^{\circ}} \\
& \text { with } \Psi(\mathbf{s})=\prod_{e_{v}>0} \psi_{v}\left(-m_{v}-s_{v}\right) \prod_{e_{v}^{\prime}>0} 2^{2 n e_{v}^{\prime}} \psi_{v}\left(-s_{v}\right)
\end{aligned}
$$

where $h=k+e-e^{\prime}$ in Case SP, $h=\left(k_{v \rho}, k_{v}+e_{v}-e_{v}^{\prime}\right)_{v \in \mathbf{b}}$ in Case UT, and $\psi_{v}$ is the polynomial $\psi_{Z}$ of Theorem 12.13 for $Z$ with $\operatorname{det}_{n}(x)^{e_{v}}$ or $\operatorname{det}_{n}(x)^{e_{v}^{\prime}}$ as its highest weight vector.

Proof. Take $\left(\tau_{n}, \varphi\right)$ as $(\alpha, \zeta)$ in Lemma 13.9. Focusing our attention on one $v \in \mathbf{a}$, we drop the subscript $v$. For $\mathfrak{z}=\operatorname{diag}[z, w]$ we have $\lambda_{\alpha}(\mathfrak{z})=$ $\left[\begin{array}{cc}1 & { }^{t} w \\ { }^{t} z & 1\end{array}\right], \mu_{\alpha}(\mathfrak{z})=\left[\begin{array}{cc}1 & w \\ z & 1\end{array}\right]$, and $\mu_{\alpha}(\mathfrak{z})^{-1}=\left[\begin{array}{cc}1-w z & 0 \\ 0 & 1-z w\end{array}\right]^{-1}\left[\begin{array}{cc}1 & -w \\ -z & 1\end{array}\right]$.
Therefore

$$
\left[\xi(\mathfrak{z})^{-1} \lambda_{\alpha}(\mathfrak{z})^{*} \cdot{ }^{t} \mu_{\alpha}(\mathfrak{z})^{-1}\right]_{\ell}=-i \cdot{ }^{t}(1-w z)^{-1}
$$

and hence $\varphi\left(\xi(\mathfrak{z})^{-1} \lambda_{\alpha}(\mathfrak{z})^{*} \cdot{ }^{t} \mu_{\alpha}(\mathfrak{z})^{-1}\right)=(-i)^{n e} j\left(\tau_{n}, \mathfrak{z}\right)^{-e}$ by (24.11b). Similarly, for $e_{v}^{\prime}>0$ we have

$$
\left[{ }^{t} \lambda_{\alpha}(\mathfrak{z}) \widehat{\mu}_{\alpha}(\mathfrak{z}) \eta(\mathfrak{z})\right]_{\ell}=\left(w-w^{*}\right) \cdot{ }^{t}(1-\overline{w z})^{-1} \eta(z),
$$

and hence

$$
\varphi\left({ }^{t} \lambda_{\alpha}(\mathfrak{z}) \widehat{\mu}_{\alpha}(\mathfrak{z}) \eta(\mathfrak{z})\right)=(4 i)^{n e^{\prime}} \delta(z)^{e^{\prime}} \delta(w)^{e^{\prime}}{\overline{j\left(\tau_{n}, \mathfrak{z}\right)}}^{-e^{\prime}}=(4 i)^{n e^{\prime}}\left(\delta^{e^{\prime}} \|_{-e^{\prime}} \tau_{n}\right)^{\circ}(z, w)
$$

Therefore we obtain the desired result from the formulas of Lemma 13.9.
25.4. To prove Theorem 20.14, we first consider Case UT. Given $h$ and $\ell$ as in that theorem, put $d=\left(d_{v}\right)_{v \in \mathbf{a}}, d_{v}=h_{v}-h_{v \rho}-\ell_{v}, e=\left(e_{v}\right)_{v \in \mathbf{a}}, e_{v}=\operatorname{Max}\left(d_{v}, 0\right)$, $e^{\prime}=\left(e_{v}^{\prime}\right)_{v \in \mathbf{a}}, e_{v}^{\prime}=\operatorname{Max}\left(-d_{v}, 0\right)$, and $k=h-d$. We now consider $E_{\mathbf{A}}$ of (24.15) on $G_{\mathbf{A}}^{2 n}$ with this $k$. Changing $\mathfrak{c}$ for its suitable multiple, we may assume, without changing $\mathcal{Z}(s, \mathbf{f}, \chi)$, that (16.24b) is satisfied. Notice that $d=e-e^{\prime}$ and $k_{v}-k_{v \rho}=$ $\ell_{v}$ for every $v \in \mathbf{a}$, so that condition (16.24a) is consistent with the assumption
on $\chi$ in Theorem 20.14. We now take $H_{b, a}$ and $\mathcal{E}_{\nu}$ of $\S 24.6$ with $n=r$, and apply $D_{e, e^{\prime}}$ to them; recall that $q=b$ if $n=r$. For $\nu=n$ we obtain, from (24.21), $\mathcal{E}_{n}=\mathcal{A}(s) \sum_{\xi \in \Xi} \chi^{*}\left(\nu_{1}(\xi)\right) N_{1}(\xi)^{-u s}\left(\delta^{\mathbf{s}} \|_{k}\left(\tau_{n}(\xi \times 1)\right)\right)(\mathfrak{z})$, and hence by (25.6) and Lemma 25.3,

$$
\left(D_{e, e^{\prime}} \mathcal{E}_{n}\right)^{\circ}=\mathcal{A}(s) \Psi(\mathbf{s}) \sum_{\xi \in \Xi} \chi^{*}\left(\nu_{1}(\xi)\right) N_{1}(\xi)^{-u s}\left(\delta^{\mathbf{s}+e^{\prime}} \|_{h} \tau_{n}\right)^{\circ} \|_{h}(\xi \times 1) .
$$

Put

$$
\mathcal{G}(z, w)=\left(D_{e, e^{\prime}} H_{b, a}\right)^{\circ}(z, \eta w) j_{\eta}^{h}(w)^{-1}, \quad \mathcal{G}_{\nu}(z, w)=\left(D_{e, e^{\prime}} \mathcal{E}_{\nu}\right)^{\circ}(z, \eta w) j_{\eta}^{h}(w)^{-1}
$$

Then

$$
\delta(w)^{m^{\prime}} \mathcal{G}_{n}\left(z,-w^{*}\right)=\mathcal{A}(s) \Psi(\mathbf{s}) \sum_{\xi \in \Xi} \chi^{*}\left(\nu_{1}(\xi)\right) N_{1}(\xi)^{-u s}\left(\delta^{s^{\prime}} \|_{h} \xi\right)(z) M^{\prime}\left(\xi z,-w^{*}\right)
$$

where $m^{\prime}=\left(h_{v}+h_{v \rho}\right)_{v \in \mathbf{a}}, \mathbf{s}^{\prime}=\mathbf{s}+e^{\prime}$, and $M^{\prime}$ is defined by (24.22a) with $m^{\prime}$ and $\mathbf{s}^{\prime}$ in place of $m$ and $\mathbf{s}$.

Now $\left(D_{e, e^{\prime}} \mathcal{E}_{\nu}\right)^{\circ}=0$ if $\nu<n$ and $d \neq 0$. In view of (25.6) this follows from

$$
\begin{equation*}
\left[D_{e, e^{\prime}}\left(\delta^{\mathbf{s}} \|_{k} \tau_{\nu}\right)\right]^{\circ}=0 \text { if } \nu<n \text { and } e+e^{\prime} \neq 0 \tag{25.7}
\end{equation*}
$$

To prove this, take $g_{\nu}$ as in Lemma 24.1. Then for $\mathfrak{z}=\operatorname{diag}[z, w]$ we have

$$
\lambda\left(\tau_{\nu}, \mathfrak{z}\right)=\left[\begin{array}{cc}
1_{n} & g_{\nu} \cdot{ }^{t} w \\
g_{\nu} \cdot{ }^{t} z & 1_{n}
\end{array}\right], \quad \mu\left(\tau_{\nu}, \mathfrak{z}\right)=\left[\begin{array}{cc}
1_{n} & g_{\nu} w \\
g_{\nu} z & 1_{n}
\end{array}\right],
$$

and we can easily verify that $\left[\xi(\mathfrak{z})^{-1} \lambda\left(\tau_{\nu}, \mathfrak{z}\right)^{*} \cdot{ }^{t} \mu\left(\tau_{\nu}, \mathfrak{z}\right)^{-1}\right]_{\ell}$ has rank $<n$. Therefore, taking $\left(\tau_{\nu}, \varphi\right)$ as $(\alpha, \zeta)$ in Lemma 13.9, we find that $\left(B_{e}\left(\delta^{\delta} \|_{k} \tau_{\nu}\right)\right)^{\circ}=0$ if $e \neq 0$. Similarly $\left[{ }^{t} \lambda\left(\tau_{\nu}, \mathfrak{z}\right) \cdot{ }^{t}{\overline{\mu\left(\tau_{\nu}, \mathfrak{z}\right)}}^{-1} \eta(\mathfrak{z})\right]_{\ell}$ has rank $<n$, and so we obtain, by the same lemma, $\left(C_{e^{\prime}}\left(\delta^{\mathbf{s}} \|_{k} \tau_{\nu}\right)\right)^{\circ}=0$ if $e^{\prime} \neq 0$. This proves (25.7).

Thus $\mathcal{G}=\mathcal{G}_{n}$ if $d \neq 0$. Returning to the setting of $\S 24.8$, given $\mathbf{f} \in \mathcal{S}_{h}^{n}(C)$, we consider

$$
\int_{\mathfrak{D}} \mathcal{G}\left(z,-w^{*}\right) f_{a}(w) \delta(w)^{m^{\prime}} \mathbf{d} w \quad\left(\mathfrak{D}=\Gamma_{0} \backslash \mathcal{H}^{r}\right)
$$

with the function $f_{a}$ associated with $\mathbf{f}$ as before. Repeating the calculations in $\S \S 24.8$ and 24.9 with $\mathcal{G}$ in place of $\mathcal{F}$, we find that

$$
\begin{align*}
& \mu c_{m^{\prime}}\left(\mathbf{s}^{\prime}\right) \chi(\theta)^{n} N(\mathfrak{b r})^{-n u s} \Psi(\mathbf{s}) f_{b}^{\prime}(z, s)  \tag{25.8}\\
& \quad=\sum_{a \in \mathcal{B}} \chi_{\mathbf{h}}(\operatorname{det}(a)) \int_{\mathfrak{D}} J_{b, a}^{\prime}\left(z,-w^{*} ; s\right) f_{a}(w) \delta(w)^{m^{\prime}} \mathbf{d} w
\end{align*}
$$

for sufficiently large $\operatorname{Re}(s)$, where $J_{b, a}^{\prime}(z, w ; s)=\left(D_{e, e^{\prime}} H_{b, a}\right)^{\circ}(z, \eta w) j_{\eta}^{h}(w)^{-1}$. If $d=0$, (25.8) becomes (24.30), and so (25.8) is true for any $d$.

In Case SP with integral $h$ the matter is simpler. Given $h$ and $\ell$ as in Theorem 20.14, put $e=h-\ell$ and $k=\ell$. We have again $\left(B_{e} \mathcal{E}_{\nu}\right)^{\circ}=0$ for $\nu<n$ if $e \neq 0$, as the proof of (25.7) is valid in this case too. Then we have (25.8) with (h, s) in place of ( $m^{\prime}, \mathrm{s}^{\prime}$ ).
25.5. Let $E_{\mathbf{A}}^{N}$ denote $E_{\mathbf{A}}$ of (16.27) defined with $G^{N}$ as $G$, and $E_{t}^{N}(\mathfrak{z})$ denote, for $t \in G_{\mathbf{h}}^{N}$ and $\mathfrak{z} \in \mathcal{H}^{N}$, the function defined by formula (17.23a), which is meaningful for integral $k$ in Cases SP and UT. In Case SP, if $t \in \alpha^{-1} D^{\prime}$ with a sufficiently small open subgroup $D^{\prime}$ of $G_{\mathbf{A}}^{N}$, then we easily see that $E_{t}^{N}=E_{1}^{N} \| \alpha$. Observe that $H_{q, a}=E_{t}^{N}$ with $t=q_{1} \sigma^{-1}$.

Suppose that $\mathbf{f}$ is an eigenform; then $f_{b}^{\prime}(z, s)=\mathfrak{T}(u s, \mathbf{f}, \chi) f_{b}(z)$ as noted at the end of $\S 24.9$, and $\mathcal{Z}(u s, \mathbf{f}, \chi)=\Lambda_{\mathfrak{c}}^{2 n}(s, \chi) \mathfrak{T}(u s, \mathbf{f}, \chi)$ by (20.23), and hence from (25.8) we obtain

$$
\begin{align*}
& \mu c_{m^{\prime}}\left(\mathbf{s}^{\prime}\right) \chi(\theta)^{n} N(\mathfrak{b r})^{-n u s} \Psi(\mathbf{s}) \mathcal{Z}(u s, \mathbf{f}, \chi) f_{b}(z)  \tag{25.8a}\\
& \quad=\sum_{a \in \mathcal{B}} \chi_{\mathbf{h}}(\operatorname{det}(a)) \int_{\mathcal{D}} \Lambda_{\mathfrak{c}}^{2 n}(s, \chi) J_{b, a}^{\prime}\left(z,-w^{*} ; s\right) f_{a}(w) \delta(w)^{m^{\prime}} \mathbf{d} w
\end{align*}
$$

Now, for the reason explained in [S97, §23.12], the integral of (25.8a) defines a meromorphic function on the whole $s$-plane. To make the result more precise, we first consider Case SP, and put $Q(s)=\mathcal{G}_{k, \kappa}^{2 n}(s) \Lambda_{c}^{2 n}$ with the symbols of Theorem 16.11, taking $2 n$ in place of $n$ there. Now $Q(s) J_{b, a}^{\prime}$ is the pullback of $Q(s) D_{e, e^{\prime}} E_{t}^{2 n}$, and hence it is meromorphic on the whole $s$-plane with possible poles only in the finite set described in Theorem 16.11. Therefore we can say the same for $\mathcal{G}_{k, \kappa}^{2 n}(s)$ times the left-hand side of (25.8a). This means that $\mathcal{G}_{k, \kappa}^{2 n}(s) c_{h}(\mathbf{s}) \Psi(\mathbf{s}) \mathcal{Z}(2 s, \mathbf{f}, \chi)$ can be continued to the whole $s$-plane with possible poles in the same set. Recall that $\mathbf{s}=\left(s_{v}\right)=s \mathbf{a}-(k-i \kappa) / 2$. Therefore our task is to show that $c_{h}(\mathbf{s}) \Psi(\mathbf{s}) \mathcal{G}_{k, \kappa}^{2 n}(s)$, with $s$ replaced by $s / 2$, produces the gamma factors as stated in Theorem 20.14.

For that purpose, let us write $g \sim g^{\prime}$ for two meromorphic functions $g$ and $g^{\prime}$ on $\mathbf{C}$ if both $g / g^{\prime}$ and $g^{\prime} / g$ are entire. Employing the explicit form of $\psi_{Z}$ in Theorem 12.13, we find that

$$
\Psi(\mathbf{s}) \sim \prod_{e_{v}>0} \prod_{b=0}^{n-1}\left(s+\left(k_{v}+i \kappa_{v}-b\right) / 2\right)
$$

Then by Lemma 24.5 we have

$$
\begin{equation*}
c_{h}(\mathbf{s}) \Psi(\mathbf{s}) \sim \prod_{v \in \mathbf{a}} \Gamma_{n}^{1}\left(s+e_{v}-\lambda_{n}+\left(k_{v}+i \kappa_{v}\right) / 2\right) \Gamma_{n}^{1}\left(s+\left(k_{v}+i \kappa_{v}\right) / 2\right)^{-1} \tag{*}
\end{equation*}
$$

Take $\gamma(s, a)$ as in Theorem 16.11 with $2 n$ in place of $n$. If $a \geq n$, we easily see that $\gamma(s, a)$ times $\left(^{*}\right)$ gives the desired $\Gamma_{v}^{h, \ell}$ of Theorem 20.14. If $(n-2) / 2 \leq a<n$, we have

$$
\begin{aligned}
\gamma(s, a) & =\Gamma_{2 a+1}^{1}(s+a / 2) \Gamma(2 s-a-1) \prod_{b=a+2}^{n} \Gamma(2 s-b) \\
& \sim \Gamma_{2 a+1}^{1}(s+a / 2) \Gamma(s-(a+1) / 2) \Gamma(s-a / 2) \prod_{b=a+2}^{n} \Gamma(2 s-b) \\
& \sim \Gamma_{2 a+2}^{1}(s+a / 2) \Gamma(s-a / 2) \prod_{b=a+2}^{n} \Gamma(2 s-b)
\end{aligned}
$$

The last product times $\Gamma_{n}^{1}(s+a / 2)^{-1}$ gives $g^{n}(s, a)$ of Theorem 20.14 for $a<n$. This completes the proof of Theorem 20.14 in Case SP.

Case UT can be handled in the same manner. In fact the argument was given in [S97, §23.12]. Though the function $J_{q, a}^{\prime}$ is different from that of [S97, (23.11.3)], the necessary $\Gamma_{v}^{h, \ell}$ in Case UT in Theorem 20.14 is a special case of $\Gamma_{v}^{h, m}$ of [S97, Theorem 20.5]. We have to take $h$ and $m$ in [S97, Theorem 20.5] to be $\left(h_{v}+h_{v \rho}\right)_{v \in \mathbf{a}}$ and $d$ with the present $d$ and $h$; also $\nu, \mu, n, r_{v}$ there are $-\left(h_{v \rho}\right)_{v \in \mathbf{a}}, \ell, 2 n, n$ here; $k$ of $[\mathrm{S} 97, \S 23.10]$ is $m^{\prime}$ here. Then $c_{h}\left(\mathbf{s}^{\prime}\right) \Psi(\mathbf{s})$ of [S97, (23.11.3)] coincides with $c_{m^{\prime}}\left(\mathbf{s}^{\prime}\right) \Psi(\mathbf{s})$ of (25.8a). In this sense no new proof is necessary in Case UT.
25.6. We next prove Theorem 23.9. Take integral $k$ in Case SP. Using the notation of Theorem 16.11, multiply both sides of (24.29a) by $\mathcal{G}_{k, \kappa}^{n+r}(s)$ and observe that $\mathcal{G}_{k, \kappa}^{n+r}(s) \Lambda_{c}^{n+r}(s, \chi) J_{q, a}$ is the pullback of $\mathcal{P}(s) \| \alpha$ with $\mathcal{P}$ of Theorem 16.11 and some $\alpha \in G$. Then we find that $c_{k}(\mathbf{s}) \mathcal{G}_{k, \kappa}^{n+r}(s) \mathcal{F}_{q}(z, s ; \mathbf{f}, \chi, C)$ can be continued to the whole $s$-plane as a meromorphic function whose poles are contained in the set described in Theorem 16.11 (with $n+r$ in place of $n$ ). Employing the explicit forms of $c_{k}(\mathbf{s})$ and $\mathcal{G}_{k, \kappa}^{n+r}(s)$, we obtain Theorem 23.9 for integral $k$ in Case SP. Case UT can be treated in a similar way. In fact, Theorem 23.9 in Case UT is a special case of [S97, Theorem 20.7] which is proved in [S97, §23.13].

We have assumed $k$ to be integral in the above. If $k$ is half-integral, we can still make all the arguments in Sections 24 and 25 valid by modifying the formulas suitably, though the analogue of Lemma 24.4, as well as the analysis of $p_{\alpha}(\mathfrak{z})$, becomes more complicated. We can eventually find the analogues of (24.29), (24.29a), and (25.8) for half-integral $k$, and obtain the desired results as stated in Theorems 21.4 and 23.9. For details, the reader is referred to [S95b]; to be precise, the formulas for half-integral $k$ corresponding to (24.29) and (25.8) are given in [S95b, (7.22) and (8.4)].
25.7. Let us now prove Theorems 23.11 and 23.12. Let the notation be as in those theorems. We first consider the case of integral $k$. Observe that the assertions of Theorems $17.7,17.8$, and 17.9 are applicable to $E_{t}^{N}=H_{q, a}$ of $\S 25.5$ in an obvious sense, since $E_{t}^{N}=E_{1}^{N} \| \alpha$ as noted there. In Case UT, Theorem 17.12 is applicable. We now evaluate (24.29) and (24.29a) at $s=\mu / 2$. Applying Lemma 24.11 to $\left(H_{q, a}\right)^{\circ}$ and taking the transform of the result by $1 \times \eta$, we find that

$$
\begin{equation*}
J_{q, a}(z, w ; \mu / 2)=\sum_{i} g_{a i}(z) h_{a i}(w) \tag{25.9}
\end{equation*}
$$

with holomorphic or nearly holomorphic $g_{a i}$ and $h_{a i}$, according to the nature of ( $n, r, \mu, F, \chi$ ). Moreover, suppose $\kappa=0$; let $W=\Phi \mathrm{Q}_{\mathrm{ab}}$ in Case SP and $W=$ $K_{\chi} \Phi \mathbf{Q}_{\mathrm{ab}}$ in Case UT with the notation of Theorems 17.9 and 17.12; then
(25.10) $\pi^{-\alpha} g_{a i}$ and $h_{a i}$ are $W$-rational, where $\alpha=\sum_{v \in \mathbf{a}}\left(m_{v}-\mu\right)(n+r) / 2$.

Suppose $\mathbf{f}$ is an eigenform and $\mu \in \Lambda(r, k)$; then $\left.\mathbf{f}\right|_{\chi} \mathfrak{T}=\Lambda_{\boldsymbol{c}}^{2 r}(s, \chi)^{-1} \mathcal{Z}(u s, \mathbf{f}, \chi)$. By Theorems 20.13 and 22.13, $\mathcal{Z}(u s, \mathbf{f}, \chi) \neq 0$ at $s=\mu / 2$ if $\mu \in \Lambda(r, k) ; \Lambda_{\mathfrak{c}}^{2 r}(s, \chi)$ is finite at $s=\mu / 2$. Also, from the formula for $c_{m}(\mathbf{s})$ in Lemma 24.5 we can easily derive that

$$
\begin{equation*}
c_{m}(\mathbf{s}) \in \pi^{r \lambda_{r}[F: \mathbf{Q}]} \mathbf{Q}^{\times} \quad \text { at } \quad s=\mu / 2 \quad \text { if } \kappa=0 \tag{25.11}
\end{equation*}
$$

and $c_{m}(\mathbf{s}) \neq 0$ at $s=\mu / 2$ for any $\kappa$. Therefore from (24.29) we obtain

$$
\begin{equation*}
E_{q}(z, \mu / 2 ; \mathbf{f}, \chi, C)=\sum_{a, i}\left\langle h_{a i}^{\prime}, f_{a}\right\rangle g_{a i}(z) \tag{25.12}
\end{equation*}
$$

where $h_{a i}^{\prime}$ is a constant multiple of $\overline{h_{a i}\left(-w^{*}\right)}$. Thus $E_{q}(z, \mu / 2 ; \mathbf{f}, \chi, C)$ belongs to $\mathcal{M}_{k}^{n}$ or $\mathcal{N}_{k}^{n, p}$ (of Lemma 24.11) with some $p$, according to the nature of ( $n, r, \mu, F$, $\chi)$. The conclusion holds for an arbitrary $\mathbf{f}$ in view of Lemma 20.12 (3). The results concerning $E_{k}^{n, r}(z, \mu / 2 ; g, \Gamma)$ follow from those for $E_{q}(z, \mu / 2 ; \mathbf{f}, \chi, C)$ by Proposition 23.6. As for $\mathcal{F}_{q}(z, \mu / 2 ; \cdots)$, we employ (24.29a) instead of (24.29). Since $\Lambda_{\mathfrak{c}}^{n+r}(s, \chi) J_{q, a}$ is the pullback of $D(z, \mu / 2 ; \cdots) \| \alpha$ or $D_{t}(z, \mu / 2 ; \cdots)$, the same technique is applicable. In this way we obtain Theorems 23.11 and 23.12 from the theorems of Section 17 mentioned above. Finally, the case of half-integral $k$ can be handled in the same way, since the analogues of (24.29) and (24.29a) can be proved as explained at the end of $\S 25.6$; for details the reader is referred to [S95b].

We insert here a lemma which will be needed in Section 28.
25.8. Lemma. The notation being as in Lemma 24.11, suppose that $n=r$ and $f \in \mathcal{N}_{k}^{2 n . p}(\Psi)$. Then we have a finite sum expression $(\pi i)^{n\left|e^{\prime}-e\right|}\left(D_{e . e^{\prime}} f\right)^{\circ}(z, w)=$ $\sum_{a=1}^{t} g_{a}(z) h_{a}(w)$ with $g_{a}, h_{a} \in \mathcal{N}_{k+e-e^{\prime}}^{n . q}(\Psi)$, where $q_{v}=p_{v}+n e_{v}$ if $e_{v}>0$ and $q_{v}=\operatorname{Max}\left(0, p_{v}-n e_{v}^{\prime}\right)$ if $e_{v}=0$.

Proof. From our definition of $D_{e . e^{\prime}}$ and Theorem 14.12 (4) we see that $D_{e . e^{\prime}} f$ is a nearly holomorphic function on $\mathcal{H}^{2 n}$ of degree $q$; also, it has Fourier expansion
whose coefficients belong to $(\pi i)^{n\left|e-e^{\prime}\right|} \Psi$. We know that $\left(D_{e, e^{\prime}} f\right)^{\circ}(z, w)$ has the automorphy property of an appropriate type. Therefore we obtain our assertion by the same type of argument as in the proof of Lemma 24.11.

## 26. Near holomorphy of Eisenstein series in Case UB

26.1. In [S97, Section 12] we defined certain Eisenstein series in Case UB similar to those of Section 23. The main purpose of this section is to prove analogues of Theorems 23.11 and 23.12 for such series. To define Eisenstein series, we have to use, as we did in [S97], certain unbounded forms of the symmetric spaces instead of the bounded domain $\mathfrak{B}_{m, n}$ of (3.7). Let us now recall some basic symbols introduced in [S97, Sections 6, 10, and 12] in the unitary case. We fix a CM-type ( $K, \tau$ ), $\tau=$ $\left\{\tau_{v}\right\}_{v \in \mathbf{a}}$, as in $\S 3.5$ throughout this section.

For $\varphi=\varphi^{*} \in G L_{n}(K)$ we denote the group $U(\varphi)$ of (1.7) by $G^{\varphi}$ in conformity with the notation of [S97]. We put $V=K_{n}^{1}$, and speak of the structure $(V, \varphi)$ and its localizations $\left(V_{v}, \varphi_{v}\right)$ for $v \in \mathbf{v}$. For $v \in \mathbf{a}$ let $r_{v}$ be the dimension of maximal $\varphi_{v}$-isotropic subspaces of $V_{v}$, and put $n=2 r_{v}+t_{v}$ with $0 \leq t_{v} \in \mathbf{Z}$. Then $\varphi_{v}$ has signature $\left(r_{v}+t_{v}, r_{v}\right)$ or $\left(r_{v}, r_{v}+t_{v}\right)$. We take and fix an element $\kappa$ of $K$ such that
(26.1) $\quad \kappa^{\rho}=-\kappa$ and $i \kappa_{v} \varphi_{v}$ has signature $\left(r_{v}+t_{v}, r_{v}\right)$ for every $v \in \mathbf{a}$.

For each $v \in$ a we fix an element $\sigma_{v} \in G L_{n}(\mathbf{C})$, as we did in [S97, §10.3], so that

$$
\kappa_{v} \sigma_{v} \varphi_{v} \sigma_{v}^{*}=-i \varphi_{v}^{\prime}, \quad \varphi_{v}^{\prime}=\left[\begin{array}{ccc}
0 & 0 & -i 1_{r_{v}}  \tag{26.2}\\
0 & \theta_{v} & 0 \\
i 1_{r_{v}} & 0 & 0
\end{array}\right]
$$

with $0<\theta_{v}=\theta_{v}^{*} \in G L_{t_{v}}(\mathbf{C})$. In this book we take $\sigma_{v}$ from $G L_{n}(\overline{\mathbf{Q}})$, which is certainly feasible; then $\theta_{v} \in G L_{t_{v}}(\overline{\mathbf{Q}})$. This is necessary for our later investigation of arithmeticity problems. We then put

$$
\begin{gather*}
\mathfrak{Z}^{\varphi}=\prod_{v \in \mathbf{a}} \mathfrak{Z}_{v}^{\varphi}, \quad \mathfrak{Z}_{v}^{\varphi}=\mathfrak{Z}\left(r_{v}, \theta_{v}\right)  \tag{26.3}\\
\mathfrak{Z}(r, \theta)=\left\{\left.\left[\begin{array}{l}
x \\
y
\end{array}\right] \in \mathbf{C}_{r}^{r+t} \right\rvert\, x \in \mathbf{C}_{r}^{r}, y \in \mathbf{C}_{r}^{t}, i\left(x^{*}-x\right)>y^{*} \theta^{-1} y\right\} \tag{26.4}
\end{gather*}
$$

where $0<\theta=\theta^{*} \in G L_{t}(\mathbf{C})$. (See [S97, §6.14] if $r_{v}=0$.) We define the origin $\mathbf{i}=\mathbf{i}^{\varphi}$ of $\mathfrak{Z}^{\varphi}$ by

$$
\mathbf{i}=\mathbf{i}^{\varphi}=\left(\mathbf{i}_{v}\right)_{v \in \mathbf{a}}, \quad \mathbf{i}_{v}=\left[\begin{array}{c}
i 1_{r_{v}}  \tag{26.5}\\
0
\end{array}\right] \in \mathfrak{Z}_{v}^{\varphi}
$$

In [S97, $\S \S 6.3$ and 10.3] we defined the action of $G_{\mathbf{A}}^{\varphi}$ on $\mathcal{Z}^{\varphi}$ and also factors of automorphy $\kappa(\alpha, z)$ and $\mu(\alpha, z)$ for $\alpha \in G_{\mathbf{A}}^{\varphi}$ and $z \in \mathfrak{Z}^{\varphi}$. Strictly speaking, these are first defined for $\alpha \in \prod_{v \in \mathbf{a}} U\left(\varphi_{v}^{\prime}\right)$, and then transferred to $G_{\mathbf{A}}^{\varphi}$ via the map $\gamma \mapsto\left(\sigma_{v} \gamma_{v} \sigma_{v}^{-1}\right)_{v \in \mathbf{a}}$ for $\gamma \in G_{\mathbf{A}}^{\varphi}$.

In the present book (in conformity with what we did in Sections 3, 4, and 5) we write $\lambda(\alpha, z)$ for $\kappa(\alpha, z)$, and put $\mu_{v}(\alpha, z)=\mu\left(\sigma_{v} \alpha_{v} \sigma_{v}^{-1}, z_{v}\right), \mu_{v \rho}(\alpha, z)=$ $\lambda\left(\sigma_{v} \alpha_{v} \sigma_{v}^{-1}, z_{v}\right)$, and $r_{v \rho}=r_{v}+t_{v}$ for $v \in \mathbf{a}$; we then put $\mu(\alpha, z)=\left(\mu_{v}(\alpha, z)\right)_{v \in \mathbf{b}}$. This is an element of $\prod_{v \in \mathbf{b}} G L_{r_{v}}(\mathbf{C})$. We also put $j(\alpha, z)=j_{\alpha}(z)=\left(j_{v}(\alpha, z)\right)_{v \in \mathbf{b}}$ with $j_{v}(\alpha, z)=\operatorname{det}\left(\mu_{v}(\alpha, z)\right)$ and $j^{k}(\alpha, z)=j_{\alpha}^{k}(z)=j_{\alpha}(z)^{k}$ for $k \in \mathbf{Z}^{\mathbf{b}}$. We then define the spaces of holomorphic automorphic forms $\mathcal{M}_{k}, \mathcal{M}_{k}(\Gamma)$, and also the space of cusp forms $\mathcal{S}_{k}, \mathcal{S}_{k}(\Gamma)$. (For the definition of a cusp form, see [S97, $\S 10.5]$.) In order to emphasize $\varphi$, we denote these by $\mathcal{M}_{k}^{\varphi}, \mathcal{S}_{k}^{\varphi}(\Gamma)$, etc.

If $\varphi$ is totally definite, we understand that $\mathcal{Z}^{\varphi}$ consists of a sigle point, written also $\mathbf{i}=\mathbf{i}^{\varphi}$, and $\mathcal{M}_{k}^{\varphi}=\mathcal{S}_{k}^{\varphi}=\mathbf{C}$. For these and other conventions, see [S97, $\S \S 6.14$, 10.3, and 10.5].

Taking $\kappa \varphi$ to be $\mathcal{T}$ of Sections 3 through 5, we can let $G^{\varphi}$ act on $\mathcal{H}=$ $\prod_{v \in \mathbf{a}} \mathfrak{B}\left(r_{v}+t_{v}, r_{v}\right)$. Notice that by our choice of $\kappa$, the signature $\left(r_{v}+t_{v}, r_{v}\right)$ is exactly ( $m_{v}, n_{v}$ ) in those sections. See the following $\S$ for the map of $\mathcal{H}$ onto $\mathfrak{Z}^{\varphi}$.
26.2. Define $\xi(z), \eta(z)$, and $\delta(z)$ for $z \in \mathcal{Z}(r, \theta)$ as in [S97, (6.1.8) and (6.3.11)]. Then formulas (12.1a, b) are valid as shown in [S97, Section 6]; also, $\Xi$ can be defined by (12.4b). Since our exposition of Sections 12 and 13 is practically axiomatic, all the definitions and formulas there are valid for $G^{\varphi}$ and $\mathcal{3}^{\varphi}$, once we have the right definition of $r(z)$. In fact, for $z=\left[\begin{array}{l}x \\ y\end{array}\right] \in \mathfrak{Z}(r, \theta)$ put

$$
r(z)=-\left[\begin{array}{c}
i 1_{r}  \tag{26.6}\\
{ }^{t} \theta^{-1} \bar{y}
\end{array}\right]^{t} \eta(z)^{-1}
$$

Then $r(z)=\xi(z)^{-1}\left[\begin{array}{c}-i 1_{r} \\ 0\end{array}\right]$; moreover, formulas (13.6a, b, c, d), (13.7), and (13.8a, b, c) are all valid in the present case. (We take $T=\mathbf{C}_{r}^{r+t}$ in those formulas.) The verification is straightforward. Notice that the entries of $\eta^{-1}$, being entries of $\operatorname{ir}(z)$, belong to $\mathcal{N}^{1}(\mathfrak{Z}(r, \theta))$; the same is true for $\xi^{-1}$, as can be seen from [S97, (6.1.11)]; thus
(26.7) The entries of $\xi(z)^{-1}$ and $\eta(z)^{-1}$ belong to $\mathcal{N}^{1}(\mathcal{Z}(r, \theta))$.

Naturally, for $z \in \mathfrak{Z}^{\varphi}$ we put $r(z)=\left(r_{v}\left(z_{v}\right)\right)_{v \in \mathbf{a}}$ with $r_{v}$ denoting the function $r$ on $\mathfrak{Z}_{v}^{\varphi}$. Then we can speak of nearly holomorphic functions on $\mathfrak{Z}^{\varphi}$, and can define $\mathcal{N}_{\omega}^{d}$ for $d \in \mathbf{Z}^{\mathbf{a}}$ and a $\overline{\mathbf{Q}}$-rational representation of $\mathfrak{K}=\prod_{v \in \mathbf{b}} G L_{r_{v}}(\mathbf{C})$.

In [S97, Lemma A2.3] we defined a holomorphic bijection $\mathfrak{t}$ of $\mathfrak{B}=\mathfrak{B}(r+t, r)$ onto $\mathfrak{Z}=\mathfrak{Z}(r, \theta)$. From the equality on line 8 from the bottom in [S97, p.216] we see that $-2 i d z={ }^{t} \kappa(z) d \mathfrak{t}(z) \mu(z)$ for $z \in \mathfrak{B}$ with holomorphic functions $\kappa$ and $\mu$ on $\mathfrak{B}$ with values in $G L_{r+t}(\mathbf{C})$ and $G L_{r}(\mathbf{C})$. Therefore $-2 i((D f) \circ \mathfrak{t})(u)=$ $D(f \circ \mathfrak{t})\left({ }^{t} \kappa(z) d u \mu(z)\right)$ for $u \in T=\mathbf{C}_{r}^{r+t}$ and $f \in C^{\infty}(\mathfrak{Z})$. Let $r_{0}$ and $\delta_{0}$ denote the functions on $\mathfrak{B}$ corresponding to $r$ and $\delta$ on $\mathcal{Z}$ as above. Take $f=\log \delta$. By [S97, (A2.3.2)] we have $f \circ \mathfrak{t}=\log \delta_{0}-\log |j(z)|^{2}+c$ with a holomorphic function $j$ on $\mathfrak{B}$ and a constant $c$. Also, by (13.8b), $\left(D \log \delta_{0}\right)(u)=\operatorname{tr}\left({ }^{t} r_{0} u\right)$ and $(D \log \delta)(u)=\operatorname{tr}\left({ }^{t} r u\right)$. Combining these, we find that

$$
\begin{equation*}
r(\mathfrak{t}(z))-(i / 2) \kappa(z) r_{0}(z) \cdot{ }^{t} \mu(z) \text { is holomorphic in } z \in \mathfrak{B} . \tag{26.8}
\end{equation*}
$$

We can naturally consider the map of $\mathcal{H}=\prod_{v \in \mathbf{a}} \mathfrak{B}\left(r_{v}+t_{v}, r_{v}\right)$ onto $\mathfrak{Z}^{\varphi}$ by taking the above $\mathfrak{t}$ at each $v \in \mathbf{a}$. Denoting the map again by $\mathfrak{t}$, from (26.8) we easily see that $f \circ \mathfrak{t} \in \mathcal{N}^{d}(\mathcal{H})$ if $f \in \mathcal{N}^{d}\left(\mathfrak{Z}^{\varphi}\right)$ with $d \in \mathbf{Z}^{\mathbf{a}}$. If $f$ is an automorphic form on $3^{\varphi}$, we have to multiply $f \circ \mathfrak{t}$ by a certain factor in order to get an automorphic form on $\mathcal{H}$, but since that factor is holomorphic on $\mathcal{H}$, it does not change the degree of near holomorphy.
26.3. Now we can introduce the notion of CM-point on $\mathfrak{Z}^{\boldsymbol{\gamma}}$ in exactly the same fashion as on $\mathcal{H}$. To be precise, we consider a CM-algebra $Y$ such that $[Y: K]=n$, and take a $K$-linear ring-injection $h: Y \rightarrow K_{n}^{n}$ such that $h\left(a^{\rho}\right)=\varphi h(a)^{*} \varphi^{-1}$ for every $a \in Y$. Then $h\left(Y^{u}\right)$ is contained in $G^{\ominus}$, and has a unique common fixed point $w$ on $\mathfrak{3}^{\boldsymbol{q}}$, which we call a CM-point on $\mathfrak{Z}^{\boldsymbol{r}}$. We can also define the symbols $\mathfrak{p}(w)$
and $\mathfrak{P}_{k}(w)$ for $k \in \mathbf{Z}^{\mathbf{b}}$ (see (11.17a)). These can be obtained by transferring the corresponding objects on $\mathcal{H}$ to those on $\mathcal{3}^{\varphi}$ by the above map $\mathfrak{t}$, or equivalently, by repeating their definitions with the symbols of Section 11 replaced by the corresponding ones on $\mathcal{Z}^{\varphi}$. Then $\mathcal{M}_{k}(\overline{\mathbf{Q}})$ and $\mathcal{M}_{k}(\Gamma, \overline{\mathbf{Q}})$ are meaningful, and all the results of Section 11 are valid in the present case. Also, $\mathcal{N}_{\omega}^{d}(\overline{\mathbf{Q}})$ and $\mathcal{N}_{\omega}^{d}(\Gamma, \overline{\mathbf{Q}})$ can be defined in the same way as in $\S 14.4$, and all the results of Section 14 in Case UB can be translated to the present case. To emphasize $\varphi$, we shall often write $\mathcal{M}_{k}^{\varphi}$ and $\mathcal{N}_{\omega}^{\varphi, d}$ for $\mathcal{M}_{k}$ and $\mathcal{N}_{\omega}^{d}$. We note a simple fact:
(26.9) If $j^{k}(h(a), w)=a^{\varepsilon}$ for every $a \in Y^{u}$ with $\varepsilon \in I_{Y}$, then $\mathfrak{P}_{k}(w)=p_{Y}(\varepsilon, \Phi)$.

Here $\Phi$ and $I_{Y}$ are as in §11.3. This follows immediately from (11.3a, b), (11.4a, b), and (11.17a).
26.4. With $(V, \varphi)$ as above, put $(W, \psi)=(V, \varphi) \oplus\left(H_{q}, \eta_{q}^{\prime}\right)$ (see §1.1) with $H_{q}=K_{2 q}^{1}$ and $\eta_{q}^{\prime}=\left[\begin{array}{cc}0 & 1_{q} \\ 1_{q} & 0\end{array}\right]$. We are going to consider Eisenstein series on $G_{\mathbf{A}}^{\psi}$ relative to this decomposition of $(W, \psi)$. For our later purposes, however, it is more natural to start from an arbitrary $(W, \psi)$, and consider such a sum decomposition for various different $(V, \varphi)$ as follows. Given $(W, \psi)$, let $l(\psi)$ be the dimension of a maximal totally $\psi$-isotropic subspace of $W$. We fix a decomposition $(W, \psi)=$ $(Z, \zeta) \oplus\left(K_{2 l}^{1}, \eta_{l}^{\prime}\right)$ with an anisotropic $\zeta$, so that $l=l(\psi)$, take a standard basis $\left\{g_{i}\right\}_{i=1}^{2 l}$ of $H_{l}$ with respect to $\eta_{l}^{\prime}$, and put
(26.10) $\quad\left(V_{r}, \varphi_{r}\right)$

$$
=(Z, \zeta) \oplus\left(\sum_{i=1}^{r}\left(K g_{l-r+i}+K g_{2 l-r+i}\right), \eta_{r}^{\prime}\right), \quad \varphi_{r}=\left[\begin{array}{ccc}
0 & 0 & 1_{r} \\
0 & \zeta & 0 \\
1_{r} & 0 & 0
\end{array}\right]
$$

where $\eta_{r}^{\prime}$ is the restriction of $\eta_{l}^{\prime}$ to $\sum_{i=1}^{r}\left(K g_{l-r+i}+K g_{2 l-r+i}\right)$. Put $p=\operatorname{dim}(Z)=$ $t_{v}+2 s_{v}$ for $v \in \mathbf{a}$ with $s_{v}=l\left(\zeta_{v}\right)$. Taking $\left(\zeta, s_{v}, t_{v}\right)$ in place of $\left(\varphi, r_{v}, t_{v}\right)$ in (26.2), we choose $\kappa, \tau_{v}, \theta_{v}$ so that

$$
\kappa_{v} \tau_{v} \zeta_{v} \tau_{v}^{*}=-i \zeta_{v}^{\prime}, \quad \zeta_{v}^{\prime}=\left[\begin{array}{ccc}
0 & 0 & -i 1_{s_{v}}  \tag{26.11}\\
0 & \theta_{v} & 0 \\
i 1_{s_{v}} & 0 & 0
\end{array}\right] .
$$

Fixing an integer $r$ such that $0 \leq r \leq l$, take $\varphi_{r}$ as $\varphi$ of $\S 26.1$; then $r_{v}=r+s_{v}$. Put then

$$
\begin{gather*}
\varphi_{v}^{\prime \prime}=\left[\begin{array}{ccc}
0 & 0 & -i 1_{r} \\
0 & \zeta_{v}^{\prime} & 0 \\
i 1_{r} & 0 & 0
\end{array}\right], \quad \varepsilon_{v}=\operatorname{diag}\left[1_{r_{v}+t_{v}},\left[\begin{array}{cc}
0 & 1_{r} \\
1_{s_{v}} & 0
\end{array}\right]\right]  \tag{26.12}\\
\xi_{v}=\operatorname{diag}\left[1_{r}, \tau_{v}, \kappa_{v}^{-1} 1_{r}\right], \quad \sigma_{v}=\varepsilon_{v} \xi_{v} \tag{26.13}
\end{gather*}
$$

Then we can easily verify that $\kappa_{v} \xi_{v} \varphi_{v} \xi_{v}^{*}=-i \varphi_{v}^{\prime \prime}$, and $\varepsilon_{v} \varphi_{v}^{\prime \prime} \varepsilon_{v}^{*}$ coincides with $\varphi_{v}^{\prime}$ of (26.2), so that (26.2) holds with the present $\sigma_{v}$. Thus we let $G^{\varphi_{r}}$ act on $\mathcal{3}^{\varphi_{r}}$ as in $\S 26.1$ for $0 \leq r \leq l$; Clearly $\left(V_{0}, \varphi_{0}\right)=(Z, \zeta)$.

In $[\mathrm{S} 97, \S 12.1]$ we put $\psi=\left[\begin{array}{ccc}0 & 0 & 1_{m} \\ 0 & \varphi & 0 \\ 1_{m} & 0 & 0\end{array}\right]$, and defined the action of $G^{\psi}$ on $\mathcal{Z}^{\psi}$. Take $m=l-r$ and $\varphi=\varphi_{r}$; then $\varphi_{l}=\delta \psi \delta^{*}$ and $\delta G^{\psi} \delta^{-1}=G^{\varphi_{l}}$ with $\delta=$ $\operatorname{diag}\left[1_{l+t},\left[\begin{array}{cc}0 & 1_{l-r} \\ 1_{r} & 0\end{array}\right]\right]$. Let $\sigma_{v}^{l}$ denote $\sigma_{v}$ of (26.13) defined with $r=l$. Then we can easily verify that $\sigma_{v}^{l} \delta$ coincides with $\tau_{v}$ of [S97, (12.1.4)], and therefore the
action of $\sigma_{v}^{l}\left(\delta \beta \delta^{-1}\right)_{v}\left(\sigma_{v}^{l}\right)^{-1}$ on $\mathbf{3}^{\psi}$ for $\beta \in G^{\psi}$ is exactly that of $\beta$ defined in [S97, p.93, line 10]. Therefore we can put $(W, \psi)=\left(V_{l}, \varphi_{l}\right)$, and identify $G^{\psi}$ with $G^{\varphi_{l}}$. Then all what we did in [S97, Section 12] is applicable to the present situation.
26.5. Throughout the rest of this section we fix an element $k$ of $Z^{\mathbf{b}}$. Before considering Eisenstein series on $G_{\mathbf{A}}^{\psi}$, we prove a few basic facts on cusp forms. Fixing a positive integer $q \leq l(\psi)$, put $(V, \varphi)=\left(V_{r}, \varphi_{r}\right)$ with $r=l-q$ and write a point $\mathfrak{z}_{v}$ of $\mathcal{Z}_{v}^{\psi}$ in the form

$$
\mathfrak{z} v=\left[\begin{array}{cc}
z_{v} & u_{v 2} \\
{ }^{t} u_{v 1} & w_{v}
\end{array}\right] \quad \text { with } z_{v} \in \mathbf{C}_{q}^{q}, u_{v 1} \in \mathbf{C}_{r_{v}+t_{v}}^{q}, u_{v 2} \in \mathbf{C}_{r_{v}}^{q}, w_{v} \in \mathbf{C}_{r_{v}}^{r_{v}+t_{v}} .
$$

Write then $\mathfrak{z}=(z, u, w)$ with $z=\left(z_{v}\right)_{v \in \mathbf{a}}, u=\left(u_{v}\right)_{v \in \mathbf{a}}, w=\left(w_{v}\right)_{v \in \mathbf{a}}$, where $u_{v}=\left[\begin{array}{ll}u_{v 1} & u_{v 2}\end{array}\right]$; then $w \in \mathfrak{Z}^{\varphi}$. In particular, for $u=0$ we put

$$
\begin{equation*}
(z, 0, w)=\operatorname{diag}[z, w] \quad\left(z \in \mathcal{H}_{q}^{\mathbf{a}}, w \in \mathfrak{Z}^{\varphi}\right) \tag{26.14}
\end{equation*}
$$

though $w_{v}$ is not necessarily square.
Let $f \in \mathcal{M}_{k}^{\psi}(\Gamma)$. In [S97, (A4.4.1)] we obtained a Fourier expansion

$$
\begin{equation*}
f(z, u, w)=\sum_{h \in \Lambda} c_{h}^{q}(u, w) \mathbf{e}_{\mathbf{a}}^{q}(h z) \quad(0<q \leq l(\psi)) \tag{26.15}
\end{equation*}
$$

with a Z-lattice $\Lambda$ in $S^{q}=\left\{x \in K_{q}^{q} \mid x^{*}=x\right\}$ and functions $c_{h}^{q}(u, w)$. More precisely, given $w \in \mathfrak{Z}^{\varphi}$ and $u \in\left(\mathbf{C}_{n}^{q}\right)^{\mathbf{a}}$, the function $f$ is defined for $i\left(z_{v}^{*}-z_{v}\right)>p_{v}$ with a positive definite hermitian matrix $p_{v}$ determined by $u_{v}$ and $w_{v}$, and we have (26.15) there. Emphasizing the dependence on $f$, we put $c_{h}^{q}(u, w)=$ $c_{h}^{q}(u, w ; f)$. (If $r_{v}=0$ for every $v \in \mathbf{a}$, then $\mathfrak{Z}^{\varphi}$ consists of a single point, $u=\left(u_{v 1}\right)_{v \in \mathbf{a}}$, and (26.15) takes the form $f(z, u)=\sum_{h \in \Lambda} c_{h}^{q}(u) \mathbf{e}_{\mathbf{a}}^{q}(h z)$.)
26.6. Proposition. The notation being as above, the following assertions hold:
(1) $c_{h}^{q}(u, w) \neq 0$ only if $h_{v} \geq 0$ for every $v \in \mathbf{a}$.
(2) $f$ is a cusp form if and only if $c_{h}^{l}\left(u, w ; f \|_{k} \alpha\right)=0$ for every $\alpha \in G^{\psi}$ and every $h$ such that $\operatorname{det}(h)=0$, where $l=l(\psi)$.
(3) $c_{0}^{q}(u, w)$ does not depend on $u$, and so we can put $c_{0}^{q}(u, w)=c_{0}^{q}(w ; f)$.
(4) $f$ is a cusp form if and only if $c_{0}^{q}\left(w ; f \|_{k} \alpha\right)=0$ for every $\alpha \in G^{\psi}$ and every $q$ such that $0<q \leq l(\psi)$.
(5) $f$ is a cusp form if and only if $c_{0}^{1}\left(w ; f \|_{k} \alpha\right)=0$ for every $\alpha \in G^{\psi}$.

Proof. All these assertions except (5) were proven in [S97, Proposition A4.5]. (The case in which $F=\mathbf{Q}$ and $\operatorname{dim}(W)=2=2 q$ was excluded in (1). But in that case $S U(\psi)$ is conjugate to $S L_{2}(\mathbf{Q})$, and so assertion (1) follows from the cusp condition.) To prove (5), let $0<t<q \leq l(\psi), \mathfrak{z}=\operatorname{diag}\left[z, \mathfrak{z}^{\prime}\right]$ and $\mathfrak{z}^{\prime}=$ $\operatorname{diag}\left[z^{\prime}, w\right]$ with $z \in \mathcal{H}_{t}^{\mathrm{a}}, z^{\prime} \in \mathcal{H}_{q-t}^{\mathrm{a}}$, and $w \in \mathfrak{J}^{\varphi}$. Then, for $f \in \mathcal{M}_{k}^{\psi}$ we have $f(\mathfrak{z})=\sum_{g \in S^{t}} c_{g}^{t}\left(0, \mathfrak{z}^{\prime}\right) \mathbf{e}_{\mathbf{a}}^{t}(g z)$. On the other hand we have a Fourier expansion

$$
f(\mathfrak{z})=\sum_{g \in S^{t}} \sum_{h \in S^{q-t}} a_{g, h}(w) \mathbf{e}_{\mathbf{a}}^{t}(g z) \mathbf{e}_{\mathbf{a}}^{q-t}\left(h z^{\prime}\right)
$$

with functions $a_{g, h}$ on $\mathfrak{Z}^{\varphi}$. Then clearly $c_{0}^{q}(w ; f)=a_{0,0}(w)$ and

$$
c_{0}^{t}\left(\operatorname{diag}\left[z^{\prime}, w\right] ; f\right)=\sum_{h \in S^{q-t}} a_{0, h}(w) \mathbf{e}_{\mathrm{a}}^{q-t}\left(h z^{\prime}\right)
$$

and hence if $c_{0}^{t}\left(\mathfrak{z}^{\prime} ; f\right)=0$, then $c_{0}^{q}(w ; f)=0$ for every $q>t$. Therefore (5) follows from (4).
26.7. Define a $\operatorname{map} \varepsilon_{q}: \mathcal{H}_{q}^{\mathbf{a}} \times \mathcal{Z}^{\varphi} \rightarrow \mathcal{Z}^{\psi}$ by $\varepsilon_{q}(z, w)=\operatorname{diag}[z, w]$. We view $U\left(\eta_{q}\right) \times G^{\varphi}$ as a subgroup of $G^{\psi}$. For $\beta \in U\left(\eta_{q}\right)$ and $\gamma \in G^{\psi}$ we easily see that

$$
\begin{equation*}
(\beta \times \gamma) \varepsilon_{q}(z, w)=\varepsilon_{q}\left(\beta^{\prime} z, \gamma w\right), \quad j\left(\beta \times \gamma, \varepsilon_{q}(z, w)\right)=j\left(\beta^{\prime}, z\right) j(\gamma, w) \tag{26.16}
\end{equation*}
$$

where $\beta^{\prime}=\operatorname{diag}\left[1_{q}, \kappa^{-1} 1_{q}\right] \beta \operatorname{diag}\left[1_{q}, \kappa 1_{q}\right]$. Now, taking $u$ of (26.15) to be 0 , for $f \in \mathcal{M}_{k}^{\psi}$ we obtain an expansion of the form

$$
\begin{equation*}
f\left(\varepsilon_{q}(z, w)\right)=\sum_{h \in S^{q}} c_{h}^{q}(w ; f) \mathbf{e}_{\mathbf{a}}^{q}(h z) \quad(0<q \leq l(\psi)) \tag{26.17}
\end{equation*}
$$

with some functions $c_{h}^{q}(w ; f)$ of $w \in \mathcal{Z}^{\varphi}$, which belong to $\mathcal{M}_{k}^{\varphi}$ as can easily be seen. Clearly $c_{0}^{q}$ coincides with that of Proposition 26.6 (3).

Transferring Proposition 11.15 to $\mathcal{Z}^{\psi}$, we have $\mathcal{M}_{k}^{\psi}=\mathcal{M}_{k}^{\psi}(\overline{\mathbf{Q}}) \otimes \mathbf{C}$. Therefore we can let an element $\sigma$ of $\operatorname{Aut}(\mathbf{C} / \overline{\mathbf{Q}})$ act on $\mathcal{M}_{k}^{\psi}$ by $\left(\sum_{\nu} c_{\nu} g_{\nu}\right)^{\sigma}=\sum_{\nu} c_{\nu}^{\sigma} g_{\nu}$ for $c_{\nu} \in \mathbf{C}$ and $g_{\nu} \in \mathcal{M}_{k}^{\psi}(\overline{\mathbf{Q}})$. Similarly we can let $\sigma$ act on $\mathcal{N}_{\omega}^{p}$. By Proposition 11.13 (1), $\mathcal{M}_{k}^{\psi}(\overline{\mathbf{Q}})$ is stable under $g \mapsto g \|_{k} \gamma$ for every $\gamma \in G^{\psi}$. Therefore we have

$$
\begin{equation*}
\left(f \|_{k} \gamma\right)^{\sigma}=f^{\sigma} \|_{k} \gamma \text { for every } \gamma \in G^{\psi} \tag{26.18}
\end{equation*}
$$

It should be noted that $f^{\sigma}$ depends on $k$. In fact, $\mathcal{M}_{k}=\mathcal{M}_{l}$ if $k_{v}+k_{v \rho}=l_{v}+l_{v \rho}$ for every $v \in \mathbf{a}$, but $\mathcal{M}_{k}(\overline{\mathbf{Q}})=\mathfrak{t} \mathcal{M}_{l}(\overline{\mathbf{Q}})$ with a constant $\mathfrak{t}$ which may not be algebraic; see Theorem 11.17 (2). Therefore $f^{\sigma}$ defined with $k$ as its weight is $t / t^{\sigma}$ times $f^{\sigma}$ defined with $l$ as its weight. Thus, whenever we speak of $f^{\sigma}$ in Case UB, it should be understood that $k$ is already fixed.
26.8. Proposition. (1) We have $c_{h}^{q}\left(w ; f^{\sigma}\right)=c_{h}^{\varphi}(w ; f)^{\sigma}$ for every $f \in \mathcal{M}_{k}^{\psi}$ and every $\sigma \in \operatorname{Aut}(\mathbf{C} / \overline{\mathbf{Q}})$.
(2) If $\Gamma$ is a congruence subgroup of $G^{\psi}$, then $\mathcal{S}_{k}^{\psi}(\Gamma)^{\sigma}=\mathcal{S}_{k}^{\psi}(\Gamma)$ for every $\sigma \in \operatorname{Aut}(\mathbf{C} / \overline{\mathbf{Q}})$ and $\mathcal{S}_{k}^{\psi}(\Gamma)=\mathcal{S}_{k}^{\psi}(\Gamma, \overline{\mathbf{Q}}) \otimes_{\overline{\mathbf{Q}}} \mathbf{C}$.

Proof. Given $\Gamma$, we can take congruence subgroups $\Gamma_{1}$ of $U\left(\eta_{q}\right)$ and $\Gamma_{2}$ of $U(\varphi)$ so that $f\left(\varepsilon_{q}(z, w)\right.$ ) as a function of $z$ (resp. $w$ ) for a fixed $w$ (resp. $z$ ) belongs to $\mathcal{M}_{k}^{\eta}\left(\Gamma_{1}\right)\left(\right.$ resp. $\left.\mathcal{M}_{k}^{\varphi}\left(\Gamma_{2}\right)\right)$ for every $f \in \mathcal{M}_{k}^{\psi}(\Gamma)$. Take a $\overline{\mathbf{Q}}$-basis $P$ of $\mathcal{M}_{k}^{\eta}\left(\Gamma_{1}, \overline{\mathbf{Q}}\right)$ over $\overline{\mathbf{Q}}$. Then $P$ is a $\mathbf{C}$-basis of $\mathcal{M}_{k}^{\eta}\left(\Gamma_{1}\right)$. We have $f\left(\varepsilon_{q}(z, w)\right)=\sum_{p \in P} g_{p}(w) p(z)$ with some functions $g_{p}$ on $\mathcal{Z}^{\varphi}$. Let $\mathcal{X}$ (resp. $\mathcal{Y}$ ) be the set of all CM-points on $\mathcal{H}_{q}^{\text {a }}$ (resp. $\mathfrak{3}^{\varphi}$ ). We have

$$
\{0\}=\bigcap_{x \in \mathcal{X}}\left\{\left(c_{p}\right) \in \mathbf{C}^{P} \mid \sum_{p \in P} c_{p} p(x)=0\right\}
$$

since $\mathcal{X}$ is dense in $\mathcal{H}_{q}^{\mathrm{a}}$. Therefore we can find a finite subset $X$ of $\mathcal{X}$ such that $\# X=\# P$ and $\operatorname{det}(p(x))_{p \in P, x \in X} \neq 0$. Then, from the equations $f\left(\varepsilon_{q}(x, w)\right)=$ $\sum_{p \in P} g_{p}(w) p(x)$ for all $x \in X$ we see that $g_{p}(w) \in \mathcal{M}_{k}^{\varphi}\left(\Gamma_{2}\right)$. We can put $p(z)=$ $\sum_{h} b_{h}(p) \mathbf{e}_{\mathbf{a}}^{q}(h z)$ with $b_{h}(p) \in \overline{\mathbf{Q}}$. Then we have

$$
\begin{equation*}
c_{h}^{q}(w ; f)=\sum_{p \in P} b_{h}(p) g_{p}(w) . \tag{26.19}
\end{equation*}
$$

Now for $x \in \mathcal{X}$ and $y \in \mathcal{Y}$ we easily see that the point $\varepsilon_{q}(x, y)$ is a CMpoint of $\mathfrak{Z}^{\psi}$ and $\mathfrak{P}_{k}\left(\varepsilon_{q}(x, y)\right)=\mathfrak{P}_{k}(x) \mathfrak{P}_{k}(y)$. Therefore, if $f \in \mathcal{M}_{k}^{\psi}(\Gamma, \overline{\mathbf{Q}})$, then $\sum_{p \in P} \mathfrak{P}_{k}(y)^{-1} \cdot g_{p}(y) \mathfrak{P}_{k}(x)^{-1} p(x)=\mathfrak{P}_{k}\left(\varepsilon_{q}(x, y)\right)^{-1} f\left(\varepsilon_{q}(x, w)\right) \in \overline{\mathbf{Q}}$ for every $(x, y) \in \mathcal{X} \times \mathcal{Y}$. Taking $x$ in $X$, we see that $\mathfrak{P}_{k}(y)^{-1} g_{p}(y) \in \overline{\mathbf{Q}}$ for every $y \in \mathcal{Y}$ and every $p \in P$. Thus $g_{p} \in \mathcal{M}_{k}^{\varphi}\left(\Gamma_{2}, \overline{\mathbf{Q}}\right)$ for every $p \in P$, and hence $c_{h}^{q}(w ; f) \in$ $\mathcal{M}_{k}^{\varphi}\left(\Gamma_{2}, \overline{\mathbf{Q}}\right)$ if $f$ is $\overline{\mathbf{Q}}$-rational. Next, let $f$ be an element of $\mathcal{M}_{k}^{\psi}(\Gamma)$ that is not necessarily $\overline{\mathbf{Q}}$-rational. By Proposition 11.15 we have $f=\sum_{a \in A} a f_{a}$ with a finite subset $A$ of $\mathbf{C}$ and $f_{a} \in \mathcal{M}_{k}^{\psi}(\Gamma, \overline{\mathbf{Q}})$. Let $\sigma \in \operatorname{Aut}(\mathbf{C} / \overline{\mathbf{Q}})$. Then $c_{h}^{q}(w ; f)=$
$\sum_{a \in A} a c_{h}^{q}\left(w ; f_{a}\right)$ and $c_{h}^{q}\left(w ; f^{\sigma}\right)=\sum_{a \in A} a^{\sigma} c_{h}^{q}\left(w ; f_{a}\right)=\sum_{a \in A} a^{\sigma} c_{h}^{q}\left(w ; f_{a}\right)^{\sigma}=$ $c_{h}^{q}(w ; f)^{\sigma}$, since $c_{h}^{q}(w ; f)$ is $\overline{\mathbf{Q}}$-rational. This proves (1). Suppose $f$ is a cusp form; then by Proposition 26.6 (4), $c_{0}^{q}\left(w ; f^{\sigma} \| \gamma\right)=c_{0}^{q}\left(w ;(f \| \gamma)^{\sigma}\right)=c_{0}^{q}(w ; f \| \gamma)^{\sigma}=0$ for every $\gamma \in G^{\psi}$, which means that $f^{\sigma}$ is a cusp form. This proves the first assertion, since $\mathcal{M}_{k}^{\psi}(\Gamma)$ is stable under $\sigma$, which can be seen from Proposition 11.15. Once this is established, the second assertion can be proved by the same technique as in the proof of Theorem 10.8 (2) by means of Lemma 10.3.
26.9. With $\left\{g_{i}\right\}$ as in $\S 26.4$ put $I_{q}^{\prime}=\sum_{i=1}^{q} K g_{i}$ and $I_{q}=\sum_{i=1}^{q} K g_{l+i}$, where $q=l-r, 0 \leq r \leq l(\psi)$; put also $(V, \varphi)=\left(V_{r}, \varphi_{r}\right)$. Then the decomposition on the first line of $\S 26.4$ can be written $(W, \psi)=(V, \varphi) \oplus\left(I_{q}^{\prime}+I_{q}, \eta_{q}^{\prime}\right)$. Define a parabolic subgroup $P_{r}^{\psi}$ of $G^{\psi}$ by

$$
\begin{equation*}
P=P_{r}^{\psi}=\left\{\alpha \in G^{\psi} \mid I_{q} \alpha=I_{q}\right\} . \tag{26.20}
\end{equation*}
$$

If we represent every element of $G^{\psi}$ according to the decomposition $W=I_{q}^{\prime} \oplus V \oplus I_{q}$, then $P$ consists of the elements of the form

$$
\alpha=\left[\begin{array}{lll}
a & b & c  \tag{26.21}\\
0 & e & f \\
0 & 0 & d
\end{array}\right]
$$

where $d$ represents the restriction of $\alpha$ to $I_{q}$. We then define a homomorphism $\pi_{r}: P \rightarrow G^{\varphi}$ and a homomorphism $\lambda_{r}: P \rightarrow K^{\times}$by

$$
\begin{equation*}
\pi_{r}(\alpha)=e, \quad \lambda_{r}(\alpha)=\operatorname{det}(d) \tag{26.22}
\end{equation*}
$$

For simplicity we shall often write $\pi$ for $\pi_{r}$. If $r=l(\psi)$, then $I_{q}=\{0\}$ and $P_{r}^{\psi}=G^{\psi}=G^{\varphi}$, and we understand that $\pi_{r}(\alpha)=\alpha$ and $\lambda_{r}(\alpha)=1$ for $\alpha \in G^{\psi}$.

For a congruence subgroup $\Gamma$ of $G^{\psi}$ we put
(26.23) $\mathcal{M}_{k}^{\varphi}(\Gamma, P)=\left\{\left.f \in \mathcal{M}_{k}^{\varphi}\left|f \|_{k} \pi(\gamma)=\lambda_{r}(\gamma)^{\ell}\right| \lambda_{r}(\gamma)\right|^{-\ell} f\right.$ for every $\left.\gamma \in \Gamma \cap P\right\}$,

$$
\begin{align*}
& \mathcal{S}_{k}^{\varphi}(\Gamma, P)=\mathcal{M}_{k}^{\varphi}(\Gamma, P) \cap \mathcal{S}_{k}^{\varphi}  \tag{26.23a}\\
& \mathcal{S}_{k}^{\varphi}(\Gamma, P, \overline{\mathbf{Q}})=\mathcal{S}_{k}^{\varphi}(\Gamma, P) \cap \mathcal{M}_{k}^{\varphi}(\overline{\mathbf{Q}}) \tag{26.24}
\end{align*}
$$

where $\ell=\left(k_{v}-k_{v \rho}\right)_{v \in \mathbf{a}}$. Then

$$
\begin{equation*}
\mathcal{S}_{k}^{\varphi}(\Gamma, P)=\mathcal{S}_{k}^{\varphi}(\Gamma, P, \overline{\mathbf{Q}}) \otimes_{\overline{\mathbf{Q}}} \mathbf{C} \tag{26.25}
\end{equation*}
$$

Indeed, let $f \in \mathcal{S}_{k}^{\varphi}(\Gamma, P)$. By Proposition 26.8 we can put $f=\sum_{a \in A} a g_{a}$ with a finite subset $A$ of $\mathbf{C}$ and $g_{a} \in \mathcal{S}_{k}^{\varphi}(\overline{\mathbf{Q}})$; we may assume that $A$ is linearly independent over $\overline{\mathbf{Q}}$. Then, for every $\gamma \in \Gamma$ we have $\sum_{a} a \lambda_{r}(\gamma)^{\ell}\left|\lambda_{r}(\gamma)\right|^{-\ell} g_{a}=f \| \pi(\gamma)=$ $\sum_{a} a g_{a} \| \pi(\gamma)$. Since $g_{a} \| \pi(\gamma) \in \mathcal{S}_{k}^{\varphi}(\overline{\mathbf{Q}})$, we have $g_{a} \| \pi(\gamma)=\lambda_{r}(\gamma)^{\ell}\left|\lambda_{r}(\gamma)\right|^{-\ell} g_{a}$, that is, $g_{a} \in \mathcal{S}_{k}^{\varphi}(\Gamma, P, \overline{\mathbf{Q}})$, which proves (26.25).

We now define a holomorphic map

$$
\begin{equation*}
\wp_{\varphi}^{\psi}: \mathfrak{Z}^{\psi} \rightarrow \mathfrak{Z}^{\varphi} \tag{26.26}
\end{equation*}
$$

by $\wp_{\varphi}^{\psi}(z, u, w)=w$ for $(z, u, w) \in \mathfrak{Z}^{\psi}$ as in $\S 26.5$. (If $r_{v}=0$, then $\wp(\mathfrak{z})_{v}=\mathbf{i}_{v}$.)
We write simply $\wp$ for $\wp_{\varphi}^{\psi}$ if there is no fear of confusion.
Given $f \in \mathcal{S}_{k}^{\varphi}(\Gamma, P)$, we put, for $(z, s) \in \mathcal{Z}^{\mathfrak{k}} \times \mathbf{C}$,

$$
\begin{align*}
\delta_{s . f}(z)= & \delta(z, s ; f)=f(\wp(z))[\delta(z) / \delta(\wp(z))]^{s \mathbf{a}-m / 2}, \quad m=\left(k_{v}+k_{v \rho \rho}\right)_{v \in \mathbf{a}}  \tag{26.27}\\
& E_{k}^{\dot{c \cdot f}}(z, s ; f, \Gamma)=\sum_{\alpha \in A} \delta_{s . f} \|_{k} \alpha, \quad A=(\Gamma \cap P) \backslash \Gamma \tag{26.28}
\end{align*}
$$

where $\delta(z)=\left(\delta\left(z_{v}\right)\right)_{v \in \mathbf{a}}$ with $\delta\left(z_{v}\right)$ defined by [S97, (6.3.11)]; we understand that $\delta(\wp(z))_{v}=1$ if $\mathcal{Z}_{v}^{\varphi}$ is trivial. These are similar to (23.11) and (23.13), and in fact, the series of (26.28) was investigated in [S97, §12.3]. For $\Gamma^{\prime} \subset \Gamma$ we can easily verify that

$$
\begin{equation*}
\left[\Gamma \cap P: \Gamma^{\prime} \cap P\right] E_{k}^{\psi, \varphi}(z, s ; f, \Gamma)=\sum_{\alpha \in \Gamma^{\prime} \backslash \Gamma} E_{k}^{\psi, \varphi}\left(z, s ; f, \Gamma^{\prime}\right) \|_{k} \alpha \tag{26.28a}
\end{equation*}
$$

If $r=l(\psi)$, then $\mathcal{S}_{k}^{\varphi}(\Gamma, P)=\mathcal{S}_{k}^{\psi}(\Gamma)$ and $\delta_{s, f}=f$, and so $E_{k}^{\psi, \psi}(z, s ; f, \Gamma)=f(z)$.
26.10. Let us now recall briefly the notion of automorphic form on $G_{\mathbf{A}}^{\varphi}$, the zeta function associated to an eigenform, and also Eisenstein series of the above type formulated as functions on $G_{\mathbf{A}}^{\psi}$. We let $n$ denote the size of $\varphi$ as we did in $\S 26.1$. Given an open subgroup $D$ of $G_{\mathbf{A}}^{\varphi}$ such that $D \cap G_{\mathbf{h}}^{\varphi}$ is compact, we denote by $\mathcal{M}_{k}^{\varphi}(D)\left(\operatorname{resp} . \mathcal{S}_{k}^{\varphi}(D)\right)$ the set of all functions $\mathbf{f}: G_{\mathbf{A}}^{\varphi} \rightarrow \mathbf{C}$ satisfying the following conditions:
(26.29a) $\mathbf{f}(\alpha x w)=j_{w}^{k}(\mathbf{i})^{-1} \mathbf{f}(x) \quad$ if $\alpha \in G^{\varphi}, \quad w \in D$, and $w(\mathbf{i})=\mathbf{i}$.
(26.29b) For every $p \in G_{\mathrm{h}}^{\varphi}$ there exists an element $f_{p}$ of $\mathcal{M}_{k}^{\varphi}$ (resp. $\mathcal{S}_{k}^{\varphi}$ ) such that $\mathbf{f}(p y)=\left(f_{p} \|_{k} y\right)(\mathbf{i})$ for every $y \in G_{\mathbf{a}}^{\varphi}$.
To find a good $D$, we first assume that

$$
\begin{equation*}
\operatorname{det}(\varphi) \in \mathfrak{g}^{\times} N_{K / F}\left(K^{\times}\right) \text {if } n \text { is odd, } \tag{26.30}
\end{equation*}
$$

which is always satisfied if we change $\varphi$ for its suitable multiple by an element of $F^{\times}$. We also fix a $\mathfrak{g}$-maximal $\mathfrak{r}$-lattice $M$ in $V$ and an integral $\mathfrak{g}$-ideal $\mathfrak{c}$, and put

$$
\begin{align*}
& \widetilde{M}=\left\{x \in V \mid \varphi(x, M) \subset \mathfrak{d}(K / F)^{-1}\right\}  \tag{26.31}\\
& D^{\varphi}=\left\{\gamma \in G_{\mathbf{A}}^{\varphi} \mid M \gamma=M, \widetilde{M}_{v}\left(\gamma_{v}-1\right) \subset \mathfrak{c}_{v} M_{v} \text { for every } v \mid \mathfrak{c}\right\}  \tag{26.32}\\
& \mathfrak{X}=\left\{\xi \in G_{\mathbf{A}}^{\varphi} \mid \xi_{v} \in D^{\varphi} \cap G_{v}^{\varphi} \text { for every } v \mid \mathfrak{c}\right\} \tag{26.33}
\end{align*}
$$

where $\mathfrak{d}(K / F)$ is the different of $K$ relative to $F$. Then we can define the action of $\mathfrak{R}\left(D^{\varphi}, \mathfrak{X}\right)$ on $\mathcal{S}_{k}^{\varphi}\left(D^{\varphi}\right)$ and a formal Dirichlet series $\mathbf{f} \mid \mathfrak{T}$ for $\mathbf{f} \in \mathcal{S}_{k}^{\varphi}\left(D^{\varphi}\right)$ as in [S97, $\S \S 11.6-11.11]$. They are similar to what was done in $\S 20.3$ of the present book. To be precise, $\mathbf{f} \mid \mathfrak{T}=\sum_{\mathfrak{a}}(\mathbf{f} \mid T(\mathfrak{a}))[\mathfrak{a}]$, where $T(\mathfrak{a})$ is the sum of all different $D^{\varphi} \tau D^{\varphi}$ with $\tau \in \mathfrak{X}$ such that the ideal $\nu^{\sigma}(\tau)$ defined by [S97, (11.11.1)] coincides with $\mathfrak{a}$. We can also define $\mathcal{Z}(s, \mathbf{f}, \chi)$ for an eigenform $\mathbf{f}$ and prove a theorem similar to Theorem 20.14; see [S97, (20.4.1) and Theorem 20.5]. (Corrections: In [S97, (20.4.1)] the condition $\mathfrak{q} \nmid \mathfrak{c}$ must be replaced by $\mathfrak{q} \nmid c \mathfrak{c h}$, where $\mathfrak{h}=F \cap$ (the conductor of $\chi$ ). Also after " $\psi=\varphi$ " in [S97, p.196, line 4] insert: "Changing $\mathfrak{c}$ for its suitable multiple, we may assume, without changing $\mathcal{Z}(s, \mathbf{f}, \chi)$, that the conductor of $\chi$ divides $\mathfrak{c}$.) Furthermore, in the setting of $\S 26.9$, given $\mathbf{f} \in \mathcal{S}_{k}^{\varphi}\left(D^{\varphi}\right)$, we can define an Eisenstein series $E_{\mathbf{A}}\left(x, s ; \mathbf{f}, \chi, D^{\psi}\right)$ for $(x, s) \in G_{\mathbf{A}}^{\psi} \times \mathbf{C}$, a Hecke character $\chi$ of $K$, and a suitably chosen open subgroup $D^{\psi}$ of $G_{\mathbf{A}}^{\psi}$, and prove its meromorphic continuation. These are similar to $E_{\mathbf{A}}$ of (23.14) and Theorem 23.9. For details, see [S97, Section 12 and Theorem 20.7]; in particular, $D^{\psi}$ is given by [S97, (20.6.8)], and we assume that $\chi$ satisfies $(16.24 \mathrm{a}, \mathrm{b})$ with $\ell=\left(k_{v}-k_{v \rho}\right)_{v \in \mathbf{a}}$ and $\mathfrak{c}$ is divisible by the ideal $\mathfrak{e}$ of [S97, Lemma 20.2]. (The last assumption is [S97, (20.3.2)].)

Let $p$ be an element of $G_{\mathbf{h}}^{\psi}$ of the form $p=\operatorname{diag}[\hat{u}, b, u]$ (with respect to the decomposition $W=I_{q}^{\prime}+V+I_{q}$ of $\left.\S 26.9\right)$ with $u \in G L_{q}(K)_{\mathbf{h}}$ and $b \in G_{\mathbf{h}}^{\varphi}$, such that $p_{v}=1$ for every $v \mid \mathfrak{c}$. Then we define a function $E_{p}\left(z, s ; \mathbf{f}, \chi, D^{\psi}\right)$ of $(z, s) \in \mathcal{Z}^{\psi} \times \mathbf{C}$ by

$$
\begin{equation*}
E_{p}\left(y(\mathbf{i}), s ; \mathbf{f}, \chi, D^{\psi}\right) j_{y}^{k}(\mathbf{i})^{-1}=E_{\mathbf{A}}\left(p y, s ; \mathbf{f}, \chi, D^{\psi}\right) \text { for every } y \in G_{\mathbf{a}}^{\psi} \tag{26.34}
\end{equation*}
$$

26.11. We are going to state an analogue of Lemma 24.11. We first put

$$
\omega=\left[\begin{array}{cc}
\psi & 0  \tag{26.35}\\
0 & -\varphi
\end{array}\right], \quad \eta=\left[\begin{array}{cc}
0 & -1_{q+n} \\
1_{q+n} & 0
\end{array}\right] \quad(n=\operatorname{dim}(V)) .
$$

Then in $[S 97,(21.1 .6)]$ we took an element $S$ of $G L_{2 q+2 n}(K)$ such that $S \eta S^{*}=\kappa \omega$. For $(\beta, \gamma) \in G^{\psi} \times G^{\varphi}$ we define an element $[\beta, \gamma]_{S}$ of $G^{\eta}$ by

$$
\begin{equation*}
[\beta, \gamma]_{S}=S^{-1} \operatorname{diag}[\beta, \gamma] S \tag{26.36}
\end{equation*}
$$

In [S97, Proposition 22.2] we defined an embedding $\iota_{U}: \mathcal{Z}^{\psi} \times \mathfrak{Z}^{\varphi} \rightarrow \mathcal{H}_{q+n}^{\mathrm{a}}$ which is compatible with (26.36) in the sense that

$$
\begin{equation*}
[\beta, \gamma]_{S_{U}}(z, w)=\iota_{U}(\beta z, \gamma w) \quad\left(z \in \mathfrak{Z}^{\psi}, w \in \mathfrak{Z}^{\varphi}\right) \tag{26.37}
\end{equation*}
$$

We note that $\iota_{U}(z, w)$ is holomorphic in $z$ and antiholomorphic in $w$ (see [S97, Proposition 6.11 and (22.2.1)]. For a function $f$ on $\mathcal{H}_{q+n}^{\mathrm{a}}$ we define its pullback $f^{\circ}$ to be a function on $\mathfrak{Z}^{\psi} \times \mathfrak{Z}^{\varphi}$ given by

$$
\begin{equation*}
f^{\circ}(z, w)=\delta(w, \wp(z))^{-m} f\left(\iota_{U}(z, w)\right) \quad\left(z \in \mathfrak{Z}^{\psi}, w \in \mathfrak{Z}^{\varphi}\right) \tag{26.38}
\end{equation*}
$$

where $m=\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$ and $\delta(w, \wp(z))$ is defined by [ $\left.\mathrm{S} 97,(6.6 .8)\right]$.
To study the nature of $f^{\circ}$ at CM-points, take CM-points $z_{0}$ and $w_{0}$ on $\mathcal{Z}^{\psi}$ and $\mathfrak{Z}^{\varphi}$ obtained from maps $r: Y \rightarrow K_{n+2 q}^{n+2 q}$ and $s: Z \rightarrow K_{n}^{n}$ with the properties described in $\S 26.3$, where $Y$ and $Z$ are CM-algebras. Define $\ell: Y \oplus Z \rightarrow K_{2 n+2 q}^{2 n+2 q}$ by $\ell(b, c)=S^{-1} \operatorname{diag}[r(b), s(c)] S$ for $(b, c) \in Y \oplus Z$; put $z_{0}=\iota_{U}\left(z_{0}, w_{0}\right)$. From (26.37) we see that $z_{0}$ is the CM-point fixed by $\ell\left((Y \oplus Z)^{u}\right)$. Let $\Phi$ (resp. $\Phi^{\prime}$ and $\Xi)$ be defined for $(Y, q)$ (resp. $(Z, s)$ and $(Y \oplus Z, \ell)$ ) as in (4.40).

We are going to show that $\Xi=\Phi+\Phi^{\prime} \rho$. For that purpose we use the notation of [S97, $\S 6.10$ and the proof of Proposition 22.2]. (The factor of automorphy $\kappa$ there is now written $\lambda$, and so we use $\lambda$ instead of $\kappa$.) Put $\alpha=[\beta, \gamma]_{S}$ with $\beta \in G^{\psi}$ and $\gamma \in G^{\varphi}$; put also $\beta_{v}^{\prime}=\tau_{v} \beta_{v} \tau_{v}^{-1}, \gamma_{v}^{\prime}=\sigma_{v} \gamma_{v} \sigma_{v}^{-1}$, and $\varepsilon_{v}=\left[\beta_{v}^{\prime}, \gamma_{v}^{\prime}\right]_{R}$. In the first paragraph of that proof, we showed that $\alpha_{v}=U_{v}^{-1} \varepsilon_{v} U_{v}$. Let $\mathfrak{z}_{1}=\iota\left(z_{0}, w_{0}\right)$; then $\mathfrak{z}_{0}=U^{-1} \mathfrak{z}_{1}$. Take $\beta=r(b)$ and $\gamma=s(c)$ with $b \in Y^{u}$ and $c \in Z^{u}$. Then $\varepsilon_{\mathfrak{z}_{1}}=\mathfrak{z}_{1}$ and $\alpha \mathfrak{z}_{0}=\mathfrak{z}_{0}$. Suppressing the subsript $v$ for simplicity, we have

$$
\begin{equation*}
\lambda_{\alpha}\left(\mathfrak{z}_{0}\right)=\lambda\left(U^{-1} \varepsilon U, U^{-1} \mathfrak{z}_{1}\right)=\lambda\left(U^{-1}, \mathfrak{z}_{1}\right) \lambda_{\varepsilon}\left(\mathfrak{z}_{1}\right) \lambda\left(U, U^{-1} \mathfrak{z}_{1}\right) . \tag{26.39}
\end{equation*}
$$

Define $M(w)$ and $N(z)$ as in [S97, (6.11.4)]. By [S97, (6.11.5)],

$$
\lambda_{\varepsilon}\left(\mathfrak{z}_{1}\right)=M\left(w_{0}\right) \operatorname{diag}\left[\lambda_{\beta}\left(z_{0}\right), \overline{\mu_{\gamma}\left(w_{0}\right)}\right] M\left(w_{0}\right)^{-1}
$$

Putting $A=\lambda\left(U^{-1}, \mathfrak{z}_{1}\right) M\left(w_{0}\right)$, we obtain the first of the following two equalities:

$$
\begin{align*}
& \lambda_{\alpha}\left(\mathfrak{z}_{0}\right)=A \cdot \operatorname{diag}\left[\lambda_{\beta}\left(z_{0}\right), \overline{\mu_{\gamma}\left(w_{0}\right)}\right] A^{-1}  \tag{26.40a}\\
& \mu_{\alpha}\left(\mathfrak{z}_{0}\right)=B \cdot \operatorname{diag}\left[\overline{\lambda_{\gamma}\left(w_{0}\right)}, \mu_{\beta}\left(z_{0}\right)\right] B^{-1} . \tag{26.40b}
\end{align*}
$$

The second formula, in which we take $B=\mu\left(U^{-1}, z_{1}\right) N\left(z_{0}\right)$, can be proved similarly. In view of the definition of $\Phi$ in (4.37) and (4.40), from (26.3ya, b) we easily see that $\Xi=\Phi+\Phi^{\prime} \rho$. Also, our definition of the symbols $\mathfrak{p}_{v}$ in $\S 11.4$ shows that

$$
\begin{align*}
\mathfrak{p}_{v \rho}\left(\mathfrak{z}_{0}\right) & =A_{v} \cdot \operatorname{diag}\left[\mathfrak{p}_{v \rho}\left(z_{0}\right), \mathfrak{p}_{v}\left(w_{0}\right)\right] A_{v}^{-1},  \tag{26.41a}\\
\mathfrak{p}_{v}\left(\mathfrak{z}_{0}\right) & =B_{v} \cdot \operatorname{diag}\left[\mathfrak{p}_{v \rho}\left(w_{0}\right), \mathfrak{p}_{v}\left(z_{0}\right)\right] B_{v}^{-1} . \tag{26.41b}
\end{align*}
$$

Here we used (11.23) to find that $p_{Z}\left(\xi \rho, \Phi^{\prime} \rho\right)=p_{Z}\left(\xi, \Phi^{\prime}\right)$ for $\xi \in I_{Z}$.
Our reasoning is valid even when $G_{\mathbf{a}}^{\varphi}$ is compact, in which case $\mathcal{Z}^{\varphi}$ consists of a single point $\mathbf{i}^{\varphi}$. In fact, we can view $\mathrm{i}^{\varphi}$ as a CM-point in the following way. Take $\zeta \in G L_{n}(K)$ so that $\zeta \varphi \zeta^{*}$ is diagonal, and define $s: K^{n} \rightarrow K_{n}^{n}$ by $s\left(a_{1}, \ldots, a_{n}\right)=\zeta^{-1} \operatorname{diag}\left[a_{1}, \ldots, a_{n}\right] \zeta$; thus $Z=K^{n}$ in the present case. Then $s\left(a^{\rho}\right)=\varphi s(a)^{*} \varphi^{-1}$, so that we can take $\mathbf{i}^{\varphi}$ to be the above $w_{0}$. (If in addition $\psi=\varphi$, then we take $z_{0}=\mathbf{i}^{\varphi}$.) From (4.37), (4.40), and (3.24a, b) we see that ( $K^{n}, \Phi^{\prime}$ ) in this case is the sum of $n$ copies of $(K, \tau \rho)$, where $(K, \tau)$ is the CM-type we fixed in $\S 3.5$, and $\mathfrak{p}_{v \rho}\left(\mathbf{i}^{\varphi}\right)=p_{K}\left(\tau_{v} \rho, \tau \rho\right) 1_{n}=p_{K}\left(\tau_{v}, \tau\right) 1_{n}$ and $\mathfrak{p}_{v}\left(\mathbf{i}^{\varphi}\right)=1$. Then $(26.41 \mathrm{a}, \mathrm{b})$ are valid, and $\mathfrak{P}_{k}\left(\mathbf{i}^{\varphi}\right)=p_{K}\left(\sum_{v \in \mathbf{a}} k_{v \rho} \tau_{v}, n \tau\right)$. Therefore, according to our general principle of $\S 11.12$, we have
(26.41c) $\mathcal{M}_{k}(\overline{\mathbf{Q}})=p_{K}\left(\sum_{v \in \mathbf{a}} k_{v \rho} \tau_{v}, n \tau\right) \overline{\mathbf{Q}}$ for every $k \in \mathbf{Z}^{\mathbf{b}}$ if $G_{\mathbf{a}}^{\varphi}$ is compact.

Notice that this is consistent with Theorem 11.17.
26.12. Lemma. If $f \in \mathcal{N}_{k}^{\eta, d}$, then $f^{\circ}(z, w)$ can be written as a finite sum $f^{\circ}(z, w)=\sum_{a=1}^{e} g_{a}(z) \overline{h_{a}(w)}$ with $g_{a} \in \mathcal{N}_{k}^{\psi, d}$ and $h_{a} \in \mathcal{N}_{k}^{\varphi, d}$. Moreover, if $f$ is $\overline{\mathbf{Q}}$-rational, we have $\mathfrak{q}^{-1} f^{\circ}(z, w)=\sum_{a=1}^{e} g_{a}^{\prime}(z) \overline{h_{a}^{\prime}(w)}$ with $g_{a}^{\prime} \in \mathcal{N}_{k}^{\psi, d}(\overline{\mathbf{Q}})$ and $h_{a} \in \mathcal{N}_{k}^{\varphi, d}(\overline{\mathbf{Q}})$, where $\mathfrak{q}=p_{K}\left(\sum_{v \in \mathbf{a}}\left(k_{v}-k_{v \rho}\right) \tau_{v}, \sum_{v \in \mathbf{a}} t_{v} \tau_{v}\right)$ with $t_{v}$ as in (26.1). This holds even if $G_{\mathbf{a}}^{\varphi}$ is compact, in which case we understand that $\mathfrak{Z}^{\varphi}=\left\{\mathbf{i}^{\varphi}\right\}$ and $\mathcal{N}_{k}^{\varphi, d}(\overline{\mathbf{Q}})=p_{K}\left(\sum_{v \in \mathbf{a}} k_{v \rho} \tau_{v}, n \tau\right) \overline{\mathbf{Q}}$. In particular, if $\psi=\varphi$ and $G_{\mathbf{a}}^{\varphi}$ is compact, then $f^{\circ}\left(\mathbf{i}^{\varphi}, \mathbf{i}^{\varphi}\right) \in p_{K}\left(\sum_{v \in \mathbf{a}}\left(k_{v}+k_{v \rho}\right) \tau_{v}, n \tau\right) \overline{\mathbf{Q}}$.

Proof. Take congruence subgroups $\Gamma^{\xi}$ of $S U(\xi)$ for $\xi=\varphi, \psi, \eta$ so that $f \in$ $\mathcal{N}_{k}^{\eta, d}\left(\Gamma^{\eta}\right)$ and $\left[\Gamma^{\psi}, \Gamma^{\varphi}\right]_{S} \subset \Gamma^{\eta}$. In view of Proposition 14.10, it is sufficient to prove the last assertion for $f \in \mathcal{N}_{k}^{\eta, d}\left(\Gamma^{\eta}, \overline{\mathbf{Q}}\right)$. We first prove that $f^{\circ}(z, w)$ (resp. $\overline{f^{\circ}(z, w)}$ ) as a function of $z$ (resp. $w$ ) belongs to $\mathcal{N}_{k}^{\psi, d}\left(\Gamma^{\psi}\right)$ (resp. $\mathcal{N}_{k}^{\varphi, d}\left(\Gamma^{\varphi}\right)$ ). The main point here is near holomorphy, since the desired automorphy property follows from [S97, (22.3.3)]. As for the assertion concerning the functions of $z$, since $\delta(w, \wp(z))$ is holomorphic in $z$, it is sufficient to show that $f\left(\iota_{U}(z, w)\right)$ as a function of $z$ belongs to $\mathcal{N}^{d}\left(\mathfrak{Z}^{\psi}\right)$. Now, if $w=\gamma\left(\mathbf{i}^{\varphi}\right)$ with $\gamma \in G^{\varphi}$, then $\iota_{U}(z, w)=\alpha \iota_{U}\left(z, \mathbf{i}^{\varphi}\right)$ with $\left.\alpha=[1, \gamma]\right]_{S}$ by (26.37). Recall also that $\iota_{U}(z, w)=$ $U^{-1} \iota(z, w)$ with $U \in G_{\mathrm{a}}^{\eta}$ and $\iota(z, w)$ defined by [S97, (6.10.2), (6.14.2), (6.14.3)]. Since $f \circ \alpha \in \mathcal{N}^{d}$ for every $\alpha \in G_{\mathbf{a}}^{\eta}$, it is sufficient to show that $f\left(\iota\left(z, \mathbf{i}^{\varphi}\right)\right)$ as a function of $z$ belongs to $\mathcal{N}^{d}\left(\mathfrak{Z}^{\psi}\right)$. Now $\iota\left(z, \mathbf{i}^{\varphi}\right)$ is holomorphic in $z$, and hence the problem can be reduced to $r(\mathfrak{z})$ for $\mathfrak{z}=\iota\left(z, \mathbf{i}^{\varphi}\right)$. Focusing our attention on $\mathfrak{Z}_{v}^{\psi}$ with a fixed $v$ and suppressing the subscript $v$, from [S97, (6.10.2)] we obtain, for $\mathfrak{z}=\iota\left(z, \mathrm{i}^{\varphi}\right), r(\mathfrak{z})=\left({ }^{t} \mathfrak{z}-\overline{\mathfrak{z}}\right)^{-1}=\operatorname{diag}\left[i \xi(z), 2 i 1_{r}\right]^{-1}$, which combined with (26.7) shows that the entries of $r(\mathfrak{z})$ are nearly holomorphic in $z$ of degree 1 . This gives the desired fact concerning $z$. Similarly $\delta(w, \wp(z))$ and $\iota(z, w)$ are holomorphic in $\bar{w}$; also $r\left(\iota\left(\mathbf{i}^{\psi}, w\right)\right)=\operatorname{diag}\left[2 i 1_{q+r}, i \tau \xi(w) \tau^{-1}\right]^{-1}$ with some $\tau \in G L_{r+t}(\mathbf{Q})$. Therefore we obtain the desired near holomorphy of $\overline{f^{\circ}(z, w)}$ in $w$.

Next, let $\left\{g_{a}\right\}_{a=1}^{e}$ be a $\overline{\mathbf{Q}}$-basis of $\mathcal{N}_{k}^{\psi, d}\left(\Gamma^{\psi}, \overline{\mathbf{Q}}\right)$. For each fixed $w$ we have $f^{\circ}(z, w)=\sum_{a=1}^{e} g_{a}(z) \overline{h_{a}(w)}$ with complex numbers $h_{a}(w)$ uniquely determined by $w$ and $a$. Since $\bigcap_{z}\left\{x \in \mathbf{C}^{e} \mid \sum_{a=1}^{e} x_{a} g_{a}(z)\right\}=\{0\}$, we can find $e$ points $z_{1}, \ldots, z_{e}$ of $\mathfrak{Z}^{\psi}$ such that $\operatorname{det}\left(g_{a}\left(z_{b}\right)\right)_{a, b=1}^{e} \neq 0$. Solving the linear equations $\overline{f^{\circ}\left(z_{b}, w\right)}=\sum_{a=1}^{e} \overline{g_{a}\left(z_{b}\right)} h_{a}(w)$, we find that $h_{a} \in \mathcal{N}_{k}^{\varphi, d}$.

To prove the $\overline{\mathbf{Q}}$-rationality of the $h_{a}$, take $z_{0}, w_{0}$, and $z_{0}$ as in $\S 26.11$. From (26.41a, b) we easily see tht $\mathfrak{P}_{k}\left(\mathfrak{z}_{0}\right)=\mathfrak{P}_{k}\left(z_{0}\right) \mathfrak{P}_{k \rho}\left(w_{0}\right)$. Now

$$
\begin{aligned}
& \mathfrak{P}_{k \rho}\left(w_{0}\right) / \mathfrak{P}_{k}\left(w_{0}\right) \\
= & \prod_{v \in \mathbf{a}}\left[\operatorname{det}\left(\mathfrak{p}_{v}\left(w_{0}\right)\right) / \operatorname{det}\left(\mathfrak{p}_{v \rho}\left(w_{0}\right)\right)\right]^{k_{v \rho}-k_{v}}=p_{K}\left(\sum_{v \in \mathbf{a}}\left(k_{v}-k_{v \rho}\right) \tau_{v}, \sum_{v \in \mathbf{a}} t_{v} \tau_{v}\right)
\end{aligned}
$$

by Theorem 11.17 (1) and the remark after it. Thus $\mathfrak{P}_{k}\left(\mathfrak{z}_{0}\right)=\mathfrak{P}_{k}\left(z_{0}\right) \mathfrak{P}_{k_{0}}(w) \mathfrak{q}$ with $\mathfrak{q}$ of our proposition. By Proposition 11.19 we may assume that every period symbol is real. Therefore

$$
\begin{equation*}
\left.\overline{\delta\left(w_{0}, \wp\left(z_{0}\right)\right)}\right)^{h} \mathfrak{P}_{k}\left(\mathfrak{z}_{0}\right)^{-1} \overline{f\left(\mathfrak{z}_{0}\right)}=\sum_{a=1}^{e} \mathfrak{P}_{k}\left(z_{0}\right)^{-1} \overline{g_{a}\left(z_{0}\right)} \cdot \mathfrak{P}_{k}\left(w_{0}\right)^{-1} \mathfrak{q}^{-1} h_{a}\left(w_{0}\right) . \tag{}
\end{equation*}
$$

Since the CM-points on $\mathfrak{Z}^{\psi}$ form a dense subset of $\mathcal{Z}^{\psi}$, we can take the above $z_{a}$ to be CM-points. In $\S 26.1$ we took algebraic $\sigma_{v}$ and $\theta_{v}$, and so the entries of $z_{v}$ and $w_{v}$ are algebraic for the same reason as in the proof of Lemma 4.13. Therefore from [S97, (6.6.1) and (6.6.8)] we see that $\delta\left(w_{0}, \wp\left(z_{0}\right)\right)$ is algebraic. Since $f$ and $g_{a}$ are $\overline{\mathbf{Q}}$-rational, $\mathfrak{P}_{k}\left(\mathfrak{z}_{0}\right)^{-1} \overline{f\left(\mathfrak{z}_{0}\right)}$ and $\mathfrak{P}_{k}\left(z_{0}\right)^{-1} \overline{g_{a}\left(z_{0}\right)}$ are algebraic. Taking $z_{0}$ of $\left(^{*}\right)$ to be those $z_{a}$, we see that $\mathfrak{P}_{k}\left(w_{0}\right)^{-1} \mathfrak{q}^{-1} h_{a}\left(w_{0}\right)$ is algebraic for every $a$ and every CM-point $w_{0}$ on $\mathcal{Z}^{\varphi}$, that is, $\mathfrak{q}^{-1} h_{a}$ is $\overline{\mathbf{Q}}$-rational for every $a$, which completes the proof. Our argument is valid even for compact $G_{\mathbf{a}}^{\varphi}$ for the reason explained at the end of $\S 26.11$.
26.13. Theorem. Define $E_{p}$ by (26.34) for $\mathbf{f} \in \mathcal{S}_{k}^{\varphi}\left(D^{\varphi}\right)$; assuming $\mathbf{f}$ to be a Hecke eigenform, put

$$
\mathcal{F}_{p}\left(z, s ; \mathbf{f}, \chi, D^{\psi}\right)=E_{p}\left(z, s ; \mathbf{f}, \chi, D^{\psi}\right) \mathcal{Z}(s, \mathbf{f}, \chi) \prod_{j=n}^{q+n-1} L_{\mathfrak{c}}\left(2 s-j, \chi_{1} \theta^{j}\right)
$$

where $\theta$ and $\chi_{1}$ are as in $\S 20.11$ and $n=\operatorname{dim}(V)$; define $E_{k}^{\psi, \varphi}(z, s ; f, \Gamma)$ by (26.28); let $\nu \in \mathbf{Z}$ and $m=\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$. Suppose that $\kappa$ of (16.24a) is 0 . Then the following assertions hold:
(i) If $\operatorname{Max}(2 n+1, q+n) \leq \nu \leq m_{v}$ and $m_{v}-\nu \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$, then $E_{p}\left(z, \nu / 2 ; \mathbf{f}, \chi, D^{\psi}\right)$ and $E_{k}^{\psi \cdot \varphi}(z, \nu / 2 ; f, \Gamma)$ belong to $\mathcal{N}_{k}^{t}$, where

$$
t=\left\{\begin{array}{l}
(q+n)(m-\nu+2) / 2 \text { if } \nu=q+n+1, F=\mathbf{Q}, \text { and } \chi_{1}=\theta^{\nu} \\
(q+n)(m-\nu \mathbf{a}) / 2 \text { otherwise }
\end{array}\right.
$$

(ii) If $2 q+2 n-m_{v} \leq \nu \leq m_{v}$ and $m_{v}-\nu \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$, then $\mathcal{F}_{p}\left(z, \nu / 2 ; \mathbf{f}, \chi, D^{\psi}\right)$ belongs to $\mathcal{N}_{k}^{t}$, except when $0 \leq \nu<q+n, \mathfrak{c}=\mathfrak{g}$, and $\chi_{1}=\theta^{\nu}$, where

$$
t=\left\{\begin{array}{l}
(q+n)(m-\nu+2) / 2 \text { if } \nu=q+n+1, F=\mathbf{Q}, \text { and } \chi_{1}=\theta^{\nu} \\
(q+n)\{m-|\nu-q-n| \mathbf{a}-(q+n) \mathbf{a}\} / 2 \text { otherwise. }
\end{array}\right.
$$

(iii) Suppose $m=\mu \mathbf{a}$ with an integer $\mu$ such that $0<\mu<q+n$; put $s_{\mu}=$ $q+n-(\mu / 2)$. Then $\mathcal{F}_{p}\left(z, s ; \mathbf{f}, \chi, D^{\psi}\right)$ has at most a simple pole at $s_{\mu}$, which occurs only when $\chi_{1}=\theta^{\mu}$. Moreover, the residue is an element of $\mathcal{M}_{k}$.

Proof. Our reasoning is the same as in $\S 25.7$. For simplicity, we suppress ( $\mathbf{f}, \chi$, $D^{\psi}$ ) in the symbols $E_{p}$ and $\mathcal{F}_{p}$. In [S97, (22.6.6)] we showed, for an eigenform $\mathbf{f}$,

$$
\begin{align*}
& \mu c_{m}(\mathbf{s}) C^{\prime}(s) \mathcal{Z}(s, \mathbf{f}, \chi) E_{p}(z, s)  \tag{26.42}\\
& \quad=\Lambda_{\mathfrak{c}}^{n}(s, \chi) \sum_{a \in \mathcal{B}} \chi_{\mathbf{h}}(\operatorname{det}(a)) \int_{\mathfrak{D}_{a}}\left(H_{p . a}\right)^{\circ}(z, w ; s) f_{a}(w) \delta(w)^{m} \mathbf{d} w .
\end{align*}
$$

Here $\mathcal{B}$ is a complete set of representatives for $G^{\varphi} \backslash G_{\mathbf{A}}^{\varphi} / D^{\varphi} ; \mathbf{f} \leftrightarrow\left(f_{a}\right)_{a \in \mathcal{B}}$ in the sense of [S97, p.80] (which is similar to the notation of §20.1); $\Gamma^{a}=G^{\varphi} \cap a D^{\varphi} a^{-1}$, $\mathfrak{D}_{a}=\Gamma^{a} \backslash \mathfrak{Z}^{\varphi}$, and $\mu=\left[\Gamma^{a} \cap \mathfrak{r}^{\times}: 1\right] ; c_{m}(\mathbf{s})$ is given by [S97, (A2.9.2)] with $\mathbf{s}=s \mathbf{a}-m / 2$ and $C^{\prime}(s)$ by $[S 97,(22.6 .5)] ; \Lambda_{\mathfrak{c}}^{n}(s, \chi)$ is as in (20.20) in Case UT; $H_{p, a}(\mathfrak{z}, s)=E_{q_{0}}(\mathfrak{z}, s)$ for $\mathfrak{z} \in \mathcal{H}_{q+n}^{\mathbf{a}}$ with $E_{q_{0}}$ of type (17.23a) obtained from an Eisenstein series on $G_{\mathbf{A}}^{q+n}$ of type (16.27) in Case UT. (See [S97, p.184, line 3]; $q_{0}$ denotes $q_{1} \Sigma_{\mathrm{h}}^{-1}$ there.) From (26.42) we immediately obtain

$$
\begin{align*}
& \mu c_{m}(\mathbf{s}) C^{\prime}(s) \mathcal{F}_{p}(z, s)  \tag{26.43}\\
& =\Lambda_{\mathfrak{c}}^{q+n}(s, \chi) \sum_{a \in \mathcal{B}} \chi_{\mathbf{h}}(\operatorname{det}(a)) \int_{\mathfrak{D}_{a}}\left(H_{p, a}\right)^{\circ}(z, w ; s) f_{a}(w) \delta(w)^{m} \mathbf{d} w
\end{align*}
$$

Now evaluate (26.42) at $s=\nu / 2$ with $\nu$ as in (i). By Theorem 17.12 (iv), $E_{q_{0}}(\mathfrak{z}, \nu / 2)$ belongs to $\pi^{\alpha} \mathcal{N}_{k}^{\eta, t}(\overline{\mathbf{Q}})$ with $\alpha=(n+q) \sum_{v \in \mathbf{a}}\left(m_{v}-\nu\right) / 2$ and $t$ as in (i), and so, by Lemma $26.12,\left(H_{p, a}\right)^{\circ}(z, w ; \nu / 2)=\pi^{\alpha} \mathfrak{q} \sum_{i=1}^{e} g_{a i}(z) \overline{h_{a i}(w)}$ with $g_{a i} \in \mathcal{N}_{k}^{\psi, t}(\overline{\mathbf{Q}})$ and $h_{a i} \in \mathcal{N}_{k}^{\varphi, t}(\overline{\mathbf{Q}})$. From [S97, (A2.9.2)] we see that for $s=\nu / 2$ with $\nu$ as in (I), $c_{m}(\mathbf{s}) \in \pi^{d_{0}} \overline{\mathbf{Q}}^{\times}$, where $d_{0}$ is the complex dimension of $\mathfrak{Z}^{\varphi}$. Since $\chi_{\mathbf{a}}(x)=x_{\mathbf{a}}^{m}\left|x_{\mathbf{a}}\right|^{-m}$ for $x \in F_{\mathbf{a}}^{\times}$, Lemma 17.5 (2) shows that $\Lambda_{\mathbf{c}}^{n}(\nu / 2, \chi) \in \pi^{\gamma} \overline{\mathbf{Q}}^{\times}$ with $\gamma=d n \nu-d n(n-1) / 2$; also $C^{\prime}(\nu / 2) \in \overline{\mathbf{Q}}^{\times}$. Thus

$$
\begin{equation*}
\mathcal{Z}(\nu / 2, \mathbf{f}, \chi) E_{p}(z, \nu / 2)=\pi^{\alpha+\gamma-d_{0}} \mathfrak{q} \sum_{a, i} \operatorname{vol}\left(\mathcal{D}_{a}\right)\left\langle h_{a i}, f_{a}\right\rangle g_{a i}^{\prime}(z) \tag{26.44}
\end{equation*}
$$

with some $g_{a i}^{\prime} \in \mathcal{N}_{k}^{\psi, t}(\overline{\mathbf{Q}})$. By [S97, Proposition 20.4 (3)], $\mathcal{Z}(s, \mathbf{f}, \chi)$ is finite and nonzero for $\operatorname{Re}(s)>n$. Therefore from (26.44) we obtain (i) for $E_{p}(z, \nu / 2)$ when $\mathbf{f}$ is a Hecke eigenform. The result holds also for an arbitrary $\mathbf{f} \in \mathcal{S}_{k}^{\varphi}\left(D^{\varphi}\right)$, since the last space is spanned by eigenforms as shown in [S97, Proposition 20.4 (1)]. This proves (i) for $E_{p}(z, \nu / 2)$, which combined with [S97, Proposition 20.10] proves (i) for $E_{k}^{\psi, \varphi}(z, \nu / 2 ; f, \Gamma)$. To prove (ii), we employ (26.43) and $D_{q_{0}}$ instead of (26.42) and $E_{q_{0}}$, where $D_{q_{0}}$ is defined by (17.24). Then we obtain (ii) from Theorem 17.12 (v). Similarly (iii) follows from Theorem 17.8 , since $\Lambda_{c}^{q+n}(s, \chi)$ is finite and nonzero at $s_{\mu}$.

Remark. For the proof of (i) we employed the nonvanishing of $\mathcal{Z}(s, \mathbf{f}, \chi)$ for $\operatorname{Re}(s)>n$. In fact, it is plausible that such nonvanishing holds for $\operatorname{Re}(s)>3 n / 4$. If that is so, we can replace $\operatorname{Max}(2 n+1, q+n)$ by $\operatorname{Max}((3 n / 2)+1, q+n)$.

We can naturally ask whether the functions of Theorem 26.13 are $\overline{\mathbf{Q}}$-rational. The answers will be given in Theorems 27.16, 29.6, and 29.7 below.
26.14. Lemma. The notation being as in §26.10, let $\mathbf{f}$ be a nonzero element of $\mathcal{S}_{k}^{\varphi}\left(D^{\varphi}\right)$ such that $\mathbf{f} \mid T(\mathfrak{a})=\lambda(\mathfrak{a}) \mathbf{f}$ with $\lambda(\mathfrak{a}) \in \mathbf{C}$ for every $\mathfrak{a}$. Then the eigenvalues $\lambda(\mathfrak{a})$ generate an algebraic number field stable under complex conjugation.

Proof. Let $\mathcal{S}_{k}^{\varphi}\left(D^{\varphi}, \overline{\mathbf{Q}}\right)$ denote the set of all $\overline{\mathbf{Q}}$-rational elements of $\mathcal{S}_{k}^{\varphi}\left(D^{\varphi}\right)$. (We call $\mathbf{f}$ as in $(26.29 \mathrm{a}, \mathrm{b}) \overline{\mathbf{Q}}$-rational if $f_{p}$ is $\overline{\mathbf{Q}}$-rational for every $p$.) By Proposition 11.15 we have $\mathcal{S}_{k}^{\varphi}\left(D^{\varphi}\right)=\mathcal{S}_{k}^{\varphi}\left(D^{\varphi}, \overline{\mathbf{Q}}\right) \otimes_{\overline{\mathbf{Q}}} \mathbf{C}$. From Proposition 11.13 and [S97, (11.9.1)] we see that $\mathcal{S}_{k}^{\varphi}\left(D^{\varphi}, \overline{\mathbf{Q}}\right)$ is stable under $\mathfrak{R}\left(D^{\varphi}, \mathfrak{X}\right)$. Since $\mathcal{S}_{k}^{\varphi}\left(D^{\varphi}, \overline{\mathbf{Q}}\right)$ is finite-dimensional over $\overline{\mathbf{Q}}$, each $\lambda(\mathfrak{a})$ is algebraic. Now, if $T(\mathfrak{a})=\sum_{\tau} D^{\varphi} \tau D^{\varphi}$, then we easily see that $T\left(\mathfrak{a}^{\rho}\right)=\sum_{\tau} D^{\varphi} \tau^{-1} D^{\varphi}$. By [S97, Proposition 11.7] we have $\langle\mathbf{f} \mid T(\mathfrak{a}), \mathbf{f}\rangle=\left\langle\mathbf{f}, \mathbf{f} \mid T\left(\mathfrak{a}^{\rho}\right)\right\rangle$, so that $\lambda\left(\mathfrak{a}^{\rho}\right)=\overline{\lambda(\mathfrak{a})}$, which gives the desired result.

## CHAPTER VII

## ARITHMETICITY OF THE CRITICAL VALUES OF ZETA FUNCTIONS AND EISENSTEIN SERIES OF GENERAL TYPES

## 27. The spaces of holomorphic Eisenstein series

27.1. Let $(W, \psi),(Z, \zeta),\left(V_{r}, \varphi_{r}\right)$, and $l=l(\psi)$ be as in $\S 26.4$; let $k \in \mathbf{Z}^{\mathbf{b}}$. Our next task is the construction of a certain subspace of $\mathcal{M}_{k}^{\psi}$, spanned by holomorphic Eisenstein series, such that $\mathcal{M}_{k}^{\psi}$ is the direct sum of $\mathcal{S}_{k}^{\psi}$ and that subspace. Put $m=\left(k_{v \rho}+k_{v}\right)_{v \in \mathbf{a}}$ and $\ell=\left(k_{v}-k_{v \rho}\right)_{v \in \mathbf{a}}$. By [S97, Proposition 10.6 (3)], $\mathcal{M}_{k}^{\psi}=$ $\mathcal{S}_{k}^{\psi}$ either if $\psi$ is anisotropic, or if $F \neq \mathbf{Q}$ and $m_{v} \neq m_{v^{\prime}}$ for some $v, v^{\prime} \in \mathbf{a}$. Therefore, our problem is meaningful only when $\psi$ is not anisotropic, that is, $l(\psi)>0$, and $m=\mu \mathbf{a}$ with $\mu \in \mathbf{Z}$, which we assume throughout this section. We present our results not only in Case UB, but also in Cases SP and UT, but for the most part give the proof in Case UB; in the other two cases we only need minor (or rather, obvious) modifications and the following changes of symbols: $G^{\psi}, G^{\varphi_{r}}, P_{r}^{\psi}, \mathfrak{Z}^{\psi}, \pi_{r}, \lambda_{r}, \wp_{r}, l, \mathcal{M}_{k}^{\psi}, \mathcal{M}_{k}^{\varphi}, \mathcal{S}_{k}^{\psi}, \mathcal{S}_{k}^{\varphi}, E_{k}^{\psi, \varphi}$, and $\operatorname{diag}[z, w]$ should be replaced by $G^{n}, G^{r}, P^{n, r}, \mathcal{H}^{n}, \pi_{r}, \lambda_{r}, \wp_{r}, n, \mathcal{M}_{k}^{n}, \mathcal{M}_{k}^{r}, \mathcal{S}_{k}^{n}, \mathcal{S}_{k}^{r}, E_{k}^{n, r}$ of $\S \S 23.1$ and 23.2, and $\operatorname{diag}[w, z] ; Z=\{0\}$ in Case UT. The case of half-integral $k$ can be included if we take $\mathcal{G}^{n}$ and $\mathcal{P}^{n, r}$ in place of $G^{n}$ and $P^{n, r}$; of course $\mu-1 / 2 \in \mathbf{Z}$ in that case. We put $[k]=k$ if $k$ is integral and $[k]=k-\mathbf{a} / 2$ otherwise; also we put $m=k$ and $\ell=[k]$ in Case SP.

We denote by $\Gamma$ an unspecified congruence subgroup of $G^{\psi}$, and by $\rho$ a real variable on $(0, \infty)$. We put $\mathbf{i}_{q}=i 1_{q}$ and view it as as an element of $\mathcal{H}_{q}^{\text {a }}$; also we write $\wp_{r}^{\psi}$ or $\wp_{r}$ for $\wp_{\varphi_{r}}^{\psi}$ (see $\S 23.1$ and (26.26)). We note here two basic formulas (see (23.4), (23.8), and [S97, (6.9.1), (12.3.4)]):

$$
\begin{align*}
& \wp_{r}(\alpha \mathfrak{z})=\pi_{r}(\alpha) \wp_{r}(\mathfrak{z})  \tag{27.1}\\
j_{\alpha}^{k}(\mathfrak{z})= & \lambda_{r}(\alpha)^{[k]}\left|\lambda_{r}(\alpha)\right|^{k-[k]} j^{k}\left(\pi_{r}(\alpha), \wp_{r}(\mathfrak{z})\right) \quad\left(\alpha \in P_{r}^{\psi}, \mathfrak{z} \in \mathfrak{Z}^{\psi}\right) . \tag{27.2}
\end{align*}
$$

The factor $\left|\lambda_{r}(\alpha)\right|^{k-[k]}$ can be eliminated if $k$ is integral. Since it is cumbersome to have such a factor in each case, we hereafter state our formulas only for integral $k$. The corresponding formulas for half-integral $k$ can be found in [S95a, Section 8]. We note that

$$
\begin{equation*}
a^{-k}|a|^{m}=a^{-\ell}|a|^{\ell} \quad \text { for every } \quad a \in K^{\times} \quad \text { if } k \in \mathbf{Z}^{\mathbf{b}} . \tag{27.3}
\end{equation*}
$$

Given a function $f: \mathcal{Z}^{\psi} \rightarrow \mathbf{C}$, we define $\Phi f: \mathcal{Z}^{\varphi_{l-1}} \rightarrow \mathbf{C}$ by

$$
(\Phi f)(w)= \begin{cases}\lim _{\rho \rightarrow \infty} f\left(\operatorname{diag}\left[\mathbf{i}_{1}, w\right]\right) & \left(\text { Case UB, } w \in \mathfrak{Z}^{\varphi_{l-1}}\right)  \tag{27.4}\\ \lim _{\rho \rightarrow \infty} f\left(\operatorname{diag}\left[w, \rho \mathbf{i}_{1}\right]\right) & \text { (Cases SP, UT } \left., w \in \mathcal{H}^{n-1}\right)\end{cases}
$$

whenever the limit exists. If $\varphi_{l-1}$ is totally definite or $\psi=\eta_{1}^{\prime}$, we ignore $w$, and so $\Phi f$ is a constant. We define $\Phi^{q}$ for $0 \leq q \leq l(\psi)$ by $\Phi^{q}=\Phi \Phi^{q-1}$, with the identity map as $\Phi^{0}$. Then $\Phi^{q} f$, if meaningful, is a function on $\mathfrak{Z}^{\varphi_{r}}, r=i-q$. If $\varphi_{r}$ is totally definite or $\psi=\eta_{q}^{\prime}$, then $\Phi^{q} f$ is a constant.
27.2. Lemma. For $f \in \mathcal{M}_{k}^{\psi}(\Gamma)$, the following assertions hold:
(1) $\left(\Phi^{q} f\right)(w)=c_{0}^{q}(w ; f)=\lim _{\rho \rightarrow \infty} f\left(\operatorname{diag}\left[\rho \mathbf{i}_{q}, w\right]\right)$.
(2) $f$ is a cusp form if and only if $\Phi\left(f \|_{k} \alpha\right)=0$ for every $\alpha \in G^{\psi}$.
(3) $\Phi^{q}\left(f \|_{k} \alpha\right)=\lambda_{r}(\alpha)^{-k}\left(\Phi^{q} f\right) \|_{k} \pi_{r}(\alpha)$ if $\alpha \in P_{r}^{\psi}$ and $q=l-r$.
(4) $\Phi^{l-r} f$ belongs to $\mathcal{M}_{k}^{\varphi_{r}}\left(\Gamma, P_{r}^{\psi}\right)$ of (23.9) or (26.23).
(5) In Cases $S P$ and UT if $f(z)=\sum_{h \in S^{n}} c(h) \mathbf{e}_{\mathbf{a}}^{n}(h z)$ for $z \in \mathcal{H}^{n}$ as in (5.22a), then $\Phi^{n} f=c(0)$, and $\left(\Phi^{n-r} f\right)(w)=\sum_{g \in S^{r}} c\left(\operatorname{diag}\left[g, 0_{n-r}\right]\right) \mathbf{e}_{\mathbf{a}}^{r}(g w)$ if $r>0$.

Proof. That $c_{0}^{q}(w ; f)=\lim _{\rho \rightarrow \infty} f\left(\operatorname{diag}\left[\rho \mathbf{i}_{q}, w\right]\right)$ can easily be seen from (26. 17) and Proposition 26.6 (1). Then we obtain (1) for $q=1$. Assuming (1) for $q=t$ and taking $q=t+1$ in the proof of Proposition 26.6, we obtain

$$
\left(\Phi^{t} f\right)\left(\operatorname{diag}\left[z^{\prime}, w\right]\right)=c_{0}^{t}\left(\operatorname{diag}\left[z^{\prime}, w\right] ; f\right)=\sum_{h \in F} a_{0, h}(w) \mathbf{e}_{\mathbf{a}}\left(h z^{\prime}\right)
$$

so that $\left(\Phi\left(\Phi^{t} f\right)\right)(w)=a_{0,0}(w)=c_{0}^{t+1}(w ; f)$ as shown there. This proves (1) for an arbitrary $q$ by induction. Assertion (5) can be proved in a similar and simpler way. Then (2) follows from Proposition 26.6 (5), or from the definition of a cusp form in §5.8. Next, let $\alpha \in P_{r}^{\psi}, \beta=\pi_{r}(\alpha)$ and $q=l-r$. By (27.1) and (27.2), for $\mathfrak{z}=\operatorname{diag}[z, w]$ with $z \in \mathcal{H}_{q}^{\mathbf{a}}$ and $w \in \mathcal{Z}^{\varphi r}$ we have $\alpha(\mathfrak{z})=\left(z_{1}, u_{1}, \beta w\right)$ with some $z_{1}$ and $u_{1}$, so that

$$
\begin{aligned}
\left(f \|_{k} \alpha\right)(\operatorname{diag}[z, w]) & =\lambda_{r}(\alpha)^{-k} j_{\beta}^{k}(w)^{-1} f\left(z_{1}, u_{1}, \beta w\right) \\
& =\lambda_{r}(\alpha)^{-k} j_{\beta}^{k}(w)^{-1} \sum_{h} c_{h}^{q}\left(u_{1}, \beta w ; f\right) \mathbf{e}_{\mathbf{a}}^{q}\left(h z_{1}\right)
\end{aligned}
$$

Observe that $z_{1}=(a z+b) a^{*}$ with $a \in G L_{q}(\mathbf{C})^{\mathbf{a}}$ and $b \in\left(\mathbf{C}_{q}^{q}\right)^{\mathbf{a}}$ that depend only on $\alpha$, and $u_{1}$ is independent of $z$, and hence $c_{0}^{q}(w ; f \| \alpha)=\lambda_{r}(\alpha)^{-k} j_{\beta}^{k}(w)^{-1} c_{0}^{q}(\beta w ; f)$, which together with (1) proves (3). Since $\lambda_{r}(\alpha) \in \mathfrak{r}^{\times}$for $\alpha \in P_{r}^{\psi} \cap \Gamma$ and $|u|^{m}=$ $|u|^{\mu \mathbf{a}}=1$ for $u \in \mathfrak{r}^{\times}$, (4) follows from (3) and (27.3).
27.3. Lemma. Let $X$ be a complete set of representatives for $P_{r}^{\psi} \backslash G^{\psi} / \Gamma$; let $f \in \mathcal{M}_{k}^{\psi}(\Gamma)$ and $q=l-r$. Then $\Phi^{q}\left(f \|_{k} \alpha\right)=0$ for every $\alpha \in G^{\psi}$ if and only if $\Phi^{q}\left(f \|_{k} \xi^{-1}\right)=0$ for every $\xi \in X$, in which case $\Phi^{q-1}\left(f \|_{k} \alpha\right)$ is a cusp form for every $\alpha \in G^{\psi}$.

Proof. Given $\alpha \in G^{\psi}$, we can put $\alpha=\gamma \xi^{-1} \beta$ with $\gamma \in \Gamma, \xi \in X$, and $\beta \in$ $P_{r}^{\psi}$. By Lemma $27.2(3), \Phi^{q}(f \| \alpha)=\Phi^{q}\left(f \| \xi^{-1} \beta\right)=c \Phi^{q}\left(f \| \xi^{-1}\right) \| \pi_{r}(\beta)$ with a constant $c$, from which our first assertion follows. Next, given any $\varepsilon \in G^{\varphi_{r+1}}$, we can find $\delta \in G^{\psi}$ such that $\pi_{r+1}(\delta)=\varepsilon$. Then $\Phi\left(\Phi^{q-1}(f \| \alpha) \| \varepsilon\right)=c^{\prime} \Phi\left(\Phi^{q-1}(f \| \alpha \delta)\right)=$ $c^{\prime} \Phi^{q}(f \| \alpha \delta)$ with a constant $c^{\prime}$, again by Lemma $27.2(3)$. Therefore, if $\Phi^{q}\left(f \| G^{\psi}\right)=$ 0 , then by Lemma $27.2(2), \Phi^{q-1}(f \| \alpha)$ must be a cusp form.
27.4. Lemma. Let $a \in R_{p}^{r}$ and $b \in R_{q}^{r}$ with $r \leq q \leq p$, where $R$ is a field with infinitely many elements. If $\operatorname{rank}\left[\begin{array}{ll}a & b\end{array}\right]=r$, then there exists an element $x \in R_{q}^{p}$ such that $\operatorname{rank}(a x+b)=r$. Moreover, if $R$ is $\mathbf{R}$ or $\mathbf{C}$, such an $x$ can be found in any nonempty open subset of $R_{q}^{p}$.

Proof. We may assume that $0<\operatorname{rank}(a)<r$, since our assertions are obvious otherwise. Take $u \in G L_{r}(R)$ and $v \in G L_{p}(R)$ so that $u a v=\left[\begin{array}{cc}1_{n} & 0 \\ 0 & 0\end{array}\right]$ with $0<n<r$. Changing $a$ and $b$ for uav and $u b$, we may assume that $a=\left[\begin{array}{cc}1_{n} & 0 \\ 0 & 0\end{array}\right]$. Put $b=\left[\begin{array}{l}c \\ d\end{array}\right]$ with $c \in R_{q}^{n}$ and $d \in R_{q}^{r-n}$; let $y$ be the upper $(n \times q)$-block of $x$. Then $\operatorname{rank}(d)=r-n$ and $a x+b=\left[\begin{array}{c}y+c \\ d\end{array}\right]$. Choosing a suitable $y$, we have $\operatorname{rank}(a x+b)=r$ as desired.
27.5. Lemma. For $\alpha \in G^{\psi}$ and $0 \leq r<l$ we have $\Phi\left(\left(\delta \circ \wp_{r} \circ \alpha\right)^{-\mathbf{a}}\left|j_{\alpha}\right|^{-2 \mathbf{a}}\right) \neq 0$ if and only if $\alpha \in P_{r}^{\psi} P_{l-1}^{\psi}$, where we understand that $\delta \circ \wp_{r} \circ \alpha=1$ if $r=0$.

Proof. Put $q=l-r$. Let $Z$ and $g_{i}$ be as in $\S 26.4$ and let $\left\{u_{1}, \ldots, u_{p}\right\}$ be a $K-$ basis of $Z$; express $\alpha$ by a matrix with respect to the basis $\left\{g_{1}, \ldots, g_{l}, u_{1}, \ldots, u_{p}\right.$, $\left.g_{l+1}, \ldots, g_{2 l}\right\}$; thus $\operatorname{dim}(Z)=p$ and $\operatorname{dim}(W)=2 l+p$. Write the $q \times(2 l+p)$-block of $\alpha$ whose $i$-th row is the $(l+p+i)$-th row of $\alpha$ for $1 \leq i \leq q$ in the form $\left[\begin{array}{lll}a & b & c\end{array}\right]$ with $a, c \in K_{l}^{q}$ and $b \in K_{p}^{p}$. This $q \times(2 l+p)$-matrix represents the restriction of $\alpha$ to $I_{q}=\sum_{i=1}^{q} K g_{l+i}$, so that $\alpha \in P_{l-1}^{\psi}$ if and only if its top row except the upper left entry of $c$ vanishes. Now fix one $v \in \mathbf{a}$; take $\kappa_{v}$ and $\tau_{v}$ as in (26.11) and define $\sigma_{v}$ by (26.13) by taking $r$ there to be $l$. Let $h$ be the $q \times(2 l+p)$-matrix whose $i$-th row is the $\left(l+s_{v}+t_{v}+i\right)$-th row of $\left(\sigma \alpha \sigma^{-1}\right)_{v}$. Then $h=\left[\begin{array}{llll}\kappa_{v}^{-1} a_{v} & b^{\prime} & c_{v} & b^{\prime \prime}\end{array}\right]$, where $b^{\prime}$ (resp. $b^{\prime \prime}$ ) consists of the first $s_{v}+t_{v}$ (resp. the last $s_{v}$ ) rows of $\kappa_{v}^{-1} b_{v} \tau_{v}^{-1}$. Since we are looking at matrices in $\mathcal{Z}_{v}^{\psi}$ and in the group acting on it, hereafter we suppress the subscript $v$. Recall that $\eta(\mathfrak{z})=i\left(\mathfrak{x}^{*}-\mathfrak{x}\right)-\mathfrak{y}^{*} \theta^{-1} \mathfrak{y}$ for $\mathfrak{z}=\left[\begin{array}{l}\mathfrak{x} \\ \mathfrak{y}\end{array}\right] \in \mathfrak{Z}^{\psi}$ (see [S97, (6.1.8)]). Let $\mathfrak{z}=\operatorname{diag}\left[\rho \mathbf{i}_{1}, w\right]$ with $w \in \mathfrak{Z}^{\varphi_{l-1}}$ and let $y$ be the upper left $(q \times q)$ block of $\eta(\alpha \mathfrak{z})_{v}^{-1}$. Put $a=\left[\begin{array}{ll}a_{1} & a_{2}\end{array}\right]$ and $c=\left[\begin{array}{ll}c_{1} & c_{2}\end{array}\right]$, where $a_{1}$ (resp. $\left.c_{1}\right)$ is the first column of $a$ (resp. c) and $\kappa^{-1} a_{1}=a^{\prime}$ and $\kappa^{-1} a_{2}=a^{\prime \prime}$. Then $h=\left[\begin{array}{llll}a^{\prime} & d & c_{1} & e\end{array}\right]$ with $d=\left[\begin{array}{ll}a^{\prime \prime} & b^{\prime}\end{array}\right]$ and $e=\left[\begin{array}{ll}c_{2} & b^{\prime \prime}\end{array}\right]$. Let us first assume that $3_{v}^{\varphi_{l-1}}$ is nontrivial. Observe that $\eta\left(\wp_{r}(\alpha \mathfrak{z})\right)$ is the lower right $r_{v} \times r_{v}$-block of $\eta\left(\alpha_{\mathfrak{z}}\right)$, where $r_{v}=$ $r+s_{v}$. Therefore by Lemma 1.3 (1) we have $2^{r_{v}} \delta\left(\wp_{r}(\alpha \mathfrak{z})\right)=2^{q+r_{"}} \delta(\alpha \mathfrak{z}) \operatorname{det}(y)=$ $2^{q+r_{v}}\left|j_{\alpha}(\mathfrak{z})\right|^{-2} \delta(\mathfrak{z}) \operatorname{det}(y)$, and hence $\left|j_{\alpha}(\mathfrak{z})\right|^{2} \delta\left(\wp_{r}(\alpha \mathfrak{z})\right)=2^{q} \delta(w) \rho \cdot \operatorname{det}(y)$. Observe that the upper $q$ rows of $\mu(\alpha, \mathfrak{z})$ is $\left[i \rho a^{\prime}+c_{1} \quad d w+e\right]$. Since $\eta(\mathfrak{z})=\operatorname{diag}[2 \rho, \eta(w)]$ and $\eta(\alpha \mathfrak{z})^{-1}=\mu(\alpha, \mathfrak{z}) \eta(\mathfrak{z})^{-1} \mu(\alpha, \mathfrak{z})^{*}$, we obtain, by a direct calculation, $y=A+B$ with

$$
A=\left(i \rho a^{\prime}+c_{1}\right)(2 \rho)^{-1}\left(i \rho a^{\prime}+c_{1}\right)^{*} \quad \text { and } \quad B=(d w+e) \eta(w)^{-1}(d w+e)^{*}
$$

Thus

$$
\left|j_{\alpha}(\mathfrak{z})\right|^{2} \delta\left(\wp_{r}\left(\alpha_{\mathfrak{z}}\right)\right)=\delta(w) \rho \cdot \operatorname{det}\left(\rho^{-1} c_{1} c_{1}^{*}+\rho C+D\right)
$$

with matrices $C$ and $D$ which do not involve $\rho$. Since $c_{1} c_{1}^{*}$ is of rank $\leq 1$, we easily see that the right-hand side is a polynomial of $\rho$. (We shall later show that $\left(\left(\delta \circ \wp_{r} \circ \alpha\right)\left|j_{\alpha}\right|^{2}\right)_{v}$ is a polynomial in $\rho$ even if $\mathcal{Z}_{v}^{\varphi_{l-1}}$ is trivial.) Suppose $\Phi\left(\left(\delta \circ \wp_{r} \circ \alpha\right)^{-\mathbf{a}}\left|j_{a}\right|^{-2 \mathbf{a}}\right) \neq 0$; then this polynomial (for any fixed $v \in \mathbf{a}$ ) must be a constant. Since $A$ and $B$ are nonnegative, we have $\rho \cdot \operatorname{det}[A+B] \geq \rho \cdot \operatorname{det}(B)$, and hence $\operatorname{det}(B)=0$. It follows that $\operatorname{rank}(d w+e)<q$. This is so for every $w \in \mathfrak{3}^{\varphi^{\prime-1}}$. By Lemma 27.4 we have $\operatorname{rank}\left[\begin{array}{ll}d & e\end{array}\right]<q$, and hence $\operatorname{rank}\left[\begin{array}{lll}a_{2} & b & c_{2}\end{array}\right]<q$. Thus we can find $u \in S L_{q}(K)$ such that the top row of $u\left[\begin{array}{lll}a_{2} & b & c_{2}\end{array}\right]$ is 0 . Put
$\beta=\operatorname{diag}\left[\widehat{u}, 1_{r+p}, u, 1_{r}\right.$; then $\beta \in P_{r}^{\psi}, \pi_{r}(\beta)=1$, and $j_{\beta}=1$. Change $\alpha$ for $\beta \alpha$ and observe that $\delta\left(\wp_{r}(\alpha \mathfrak{z})\right)\left|j_{\alpha}(\mathfrak{z})\right|^{2}$ does not change. Thus we may assume that the top row of $\left[\begin{array}{lll}a_{2} & b & c_{2}\end{array}\right]$ is 0 . Then the top row of $\left[\begin{array}{ll}d & e\end{array}\right]$ is 0 , and so the top row of $B$ is 0 . Let $\nu$ and $\nu^{\prime}$ be the first components of $a_{1}$ and $c_{1}$. Suppose $\nu \neq 0$. Then the top row of $A$ is 0 for $\rho=i \kappa \nu^{\prime} / \nu$. Thus $\rho \cdot \operatorname{det}(y)=0$ for such a $\rho$, so that the constant $\rho \cdot \operatorname{det}(y)$ must be 0 . On the other hand, the quantity is positive for $0<\rho \in \mathbf{R}$, which is a contradiction. Therefore $\nu=0$, and hence the top row of [ $\left.\begin{array}{lll}a & b & c_{2}\end{array}\right]$ is 0 , which means that the present $\alpha$ belongs to $P_{l-1}^{\psi}$. Thus the original $\alpha$ belongs to $P_{r}^{\psi} P_{l-1}^{\psi}$ as expected.

We have assumed that $\mathcal{Z}_{v}^{\varphi_{l-1}}$ is nontrivial. If $\mathcal{Z}_{v}^{\varphi_{l-1}}$ is trivial, then $l=1$ and $s_{v}=0$, and so $r=0$ and $q=1$. In this case $\mu(\alpha, \mathfrak{z})_{v}=\left(\rho a^{\prime} i+c_{1}\right)_{v}$, and hence $\delta\left(\wp_{r}\left(\alpha_{\mathfrak{z}}\right)\right)_{v}\left|j_{\alpha}(\mathfrak{z})_{v}\right|^{2}=\left|\left(\rho a^{\prime} i+c_{1}\right)_{v}\right|^{2}$, which is again a polynomial in $\rho$. Thus $\Phi\left(\left(\delta \circ \wp_{r} \circ \alpha\right)^{-\mathbf{a}}\left|j_{\alpha}\right|^{-2 \mathbf{a}}\right) \neq 0$ only if $a^{\prime}=0$, in which case $a=0$. Since $s_{v}=0, \zeta$ is anisotropic. Now $\left[\begin{array}{lll}a & b & c\end{array}\right]$ represents $g_{2} \alpha$, so that $b \zeta b^{*}=0$. Thus $b=0$, which shows that $\alpha \in P_{0}^{\psi}$ as expected. This proves the 'only if'-part.

Conversely, let $\alpha=\gamma \varepsilon$ with $\gamma \in P_{r}^{\psi}$ and $\varepsilon \in P_{l-1}^{\psi}$; put $\xi=\pi_{l-1}(\varepsilon)$. Then by (27.1) and (27.2), for $\mathfrak{z}=\operatorname{diag}\left[\rho \mathbf{i}_{1}, w\right]$ we have $\wp_{r}(\varepsilon \mathfrak{z})=\wp_{r}^{\varphi_{l-1}}\left(\wp_{l-1}(\varepsilon \mathfrak{z})\right)=$ $\wp_{r}^{\varphi_{l-1}}(\xi w)$ and $\delta\left(\wp_{r}\left(\alpha_{\mathfrak{z}}\right)\right)\left|j_{\alpha}(\mathfrak{z})\right|^{2}=\left|\lambda_{r}(\gamma) \lambda_{l-1}(\varepsilon) j_{\xi}(w)\right|^{2} \delta\left(\wp_{r}^{\varphi_{l-1}}(\xi w)\right)$ which proves the 'if'-part.
27.6. Lemma. Let $\delta_{f}^{\psi}(\mathfrak{z})=\delta(\mathfrak{z}, s ; f)$ with the notation of (26.27) with $f \in$ $\mathcal{S}_{k}^{\varphi}, \varphi=\varphi_{r}$, and let $\alpha \in G^{\psi}$. Then for $\operatorname{Re}(s)>0$ and $l>t \geq r$ we have $\Phi^{l-t}\left\{\delta(\mathfrak{z})^{m / 2-s \mathbf{a}} \cdot\left[\delta_{f}^{\psi} \|_{k} \alpha\right]\right\} \neq 0$ only if $\alpha=\beta \gamma$ with $\beta \in P_{r}^{\psi}$ and $\gamma \in P_{t}^{\psi}$, in which case for $\mathfrak{z}=\operatorname{diag}\left[w^{\prime}, w\right]$ with $\left(w^{\prime}, w\right) \in \mathcal{H}_{l-t}^{\mathbf{a}} \times \mathfrak{Z}^{\omega}, \omega=\varphi_{t}$, we have

$$
\begin{align*}
& \delta(\mathfrak{z})^{m / 2-s \mathbf{a}}\left[\delta_{f}^{\psi} \|_{k} \alpha\right]=\Phi^{l-t}\left\{\delta(\mathfrak{z})^{m / 2-s \mathbf{a}}\left[\delta_{f}^{\psi} \|_{k} \alpha\right]\right\}  \tag{27.5}\\
= & \delta(w)^{m / 2-s \mathbf{a}}\left|\lambda_{r}(\beta) \lambda_{t}(\gamma)\right|^{m-2 s \mathbf{s}}\left(\lambda_{r}(\beta) \lambda_{t}(\gamma)\right)^{-k} \delta\left(w, s ; f \|_{k} \pi_{r}(\beta)\right) \|_{k} \pi_{t}(\gamma)
\end{align*}
$$

Proof. Let $\alpha=\beta \gamma$ and $\mathfrak{z}=\operatorname{diag}\left[w^{\prime}, w\right]$ as above. Put $g=f \| \pi_{r}(\beta)$. By [S97, (12.3.5)] or (23.12) we have $\delta_{f}^{\psi} \| \beta=\lambda_{r}(\beta)^{-k}\left|\lambda_{r}(\beta)\right|^{m-2 s \mathbf{s}} \delta_{g}^{\psi}$, and

$$
\begin{equation*}
\delta(\mathfrak{z})^{m / 2-s \mathbf{a}} \cdot\left(\delta_{g}^{\psi} \| \gamma\right)=j_{\gamma}^{k}(\mathfrak{z})^{-1}\left|j_{\gamma}(\mathfrak{z})\right|^{m-2 s \mathbf{a}} g\left(\wp_{r}(\gamma \mathfrak{z})\right) \delta\left(\wp_{r}(\gamma \mathfrak{z})\right)^{m / 2-s \mathbf{a}} \tag{}
\end{equation*}
$$

Put $\xi=\pi_{t}(\gamma)$. Then $\wp_{r}(\gamma \mathfrak{z})=\wp_{r}^{\omega}\left(\wp_{t}(\gamma \mathfrak{z})\right)=\wp_{r}^{\omega}\left(\xi \wp_{t}(\mathfrak{z})\right)=\wp_{r}^{\omega}(\xi w)$, and hence, by (27.2), the quantity of (*) equals $\lambda_{t}(\gamma)^{-k}\left|\lambda_{t}(\gamma)\right|^{m-2 s \mathbf{a}} \delta(w)^{m / 2-s \mathbf{a}} \cdot\left(\delta_{g}^{\omega} \| \xi\right)$. From these we obtain (27.5). Now by [S97, Proposition 10.6 (1)], $\left|\delta(z)^{m / 2} f(z)\right| \leq C$ for every $z \in \mathfrak{Z}^{\varphi}$ with a constant $C$. Therefore, for $\alpha \in G^{\psi}$ we have

$$
\begin{aligned}
\left|\delta(\mathfrak{z})^{m / 2-s \mathbf{a}}\left[\delta_{f}^{\psi} \| \alpha\right]\right| & =\left|f\left(\wp_{r}(\alpha \mathfrak{z})\right) \delta\left(\wp_{r}(\alpha \mathfrak{z})\right)^{m / 2-s \mathbf{a}} j_{\alpha}(\mathfrak{z})^{-2 s \mathbf{a}}\right| \\
& \leq C\left|\delta\left(\wp_{r}(\alpha \mathfrak{z})\right)^{-s \mathbf{a}} j_{\alpha}(\mathfrak{z})^{-2 s \mathbf{a}}\right|
\end{aligned}
$$

By Lemma 27.5 the last quantity, with $\mathfrak{z}=\operatorname{diag}\left[\rho \mathbf{i}_{1}, \mathfrak{z}_{1}\right]$ and $\operatorname{Re}(s)>0$, tends to 0 as $\rho \rightarrow \infty$ if $\alpha \notin P_{r}^{\psi} P_{l-1}^{\psi}$. This proves the case $t=l-1$. Suppose that our lemma is true for $\Phi^{l-t}, t>r$, and $\Phi^{l-t+1}\left(\delta(\mathfrak{z})^{m / 2-s \mathbf{a}}\left[\delta_{f}^{\psi} \| \alpha\right]\right) \neq 0$. Then $\Phi^{l-t}\left(\delta(\mathfrak{z})^{m / 2-s \mathbf{a}}\left[\delta_{f}^{\psi} \| \alpha\right]\right) \neq 0$, and so $\alpha=\beta \gamma$ with $\beta \in P_{r}^{\psi}$ and $\gamma \in P_{t}^{\psi}$; also, by (27.5), for $w \in \mathfrak{Z}^{\omega}$ we have

$$
a b^{-s} \Phi\left\{\delta(w)^{m / 2-s \mathbf{a}} \delta\left(w, s ; f \| \pi_{r}(\beta)\right) \| \pi_{t}(\gamma)\right\}=\Phi\left\{\Phi^{l-t}\left(\delta(\mathfrak{z})^{m / 2-s \mathbf{a}}\left[\delta_{f}^{\psi} \| \alpha\right]\right)\right\}
$$

$$
=\Phi^{l-t+1}\left(\delta(\mathfrak{z})^{m / 2-s \mathbf{a}}\left[\delta_{f}^{\psi} \| \alpha\right]\right) \neq 0
$$

with some $a \in \mathbf{C}^{\times}$and $0<b \in \mathbf{R}$. Applying our result in the case $\Phi^{l-t}=\Phi$ to the quantity involving $w$ on the left-hand side, we have $\pi_{t}(\gamma)=\xi^{\prime} \eta^{\prime}$ with $\xi^{\prime} \in P_{r}^{\omega}$ and $\eta^{\prime} \in P_{t-1}^{\omega}$. Now we can easily verify that

$$
\begin{gather*}
P_{t}^{\psi} \cap P_{r}^{\psi}=\left\{\alpha \in P_{t}^{\psi} \mid \pi_{t}(\alpha) \in P_{r}^{\omega}\right\},  \tag{27.6a}\\
\pi_{t}\left(P_{t}^{\psi} \cap P_{r}^{\psi}\right)=P_{r}^{\omega} \tag{27.6b}
\end{gather*}
$$

and hence we can put $\xi^{\prime}=\pi_{t}(\xi)$ and $\eta^{\prime}=\pi_{t}(\eta)$ with $\xi \in P_{t}^{\psi} \cap P_{r}^{\psi}$ and $\eta \in$ $P_{t}^{\psi} \cap P_{t-1}^{\psi}$. Then $\alpha=\beta \xi \xi^{-1} \gamma, \xi^{-1} \gamma \in P_{t}^{\psi}, \beta \xi \in P_{r}^{\psi}$, and $\pi_{t}\left(\xi^{-1} \gamma\right)=\eta^{\prime} \in P_{t-1}^{\omega}$, and hence $\xi^{-1} \gamma \in P_{t-1}^{\psi}$ by (27.6a). Therefore induction proves our lemma.
27.7. Lemma. Put $\omega=\varphi_{t}, P=P_{t}^{\psi}$, and $Q=P_{r}^{\psi}, r \leq t<l$. For each $\xi$ in the set $X$ of Lemma 27.3 let $Z_{\xi}$ be a complete set of representatives for $(P \cap Q) \backslash Q /\left(\xi \Gamma \xi^{-1} \cap Q\right)$. For every $\zeta \in Z_{\xi}$ such that $\zeta \xi \Gamma \cap P \neq \varnothing$ choose and fix an element $\eta \in \zeta \xi \Gamma \cap P$. Let $Y$ denote the set of all such $\eta$ 's. Then the following assertios hold:
(1) $Y$ is a finite set, $P=\bigsqcup_{\eta \in Y}(P \cap Q) \eta(\Gamma \cap P)$, and $G^{\omega}=\bigsqcup_{\eta \in Y} P_{t}^{\omega} \pi_{t}(\eta(\Gamma \cap P))$.
(2) If $R_{\eta}$ is a complete set of representatives for $\left(\eta \Gamma \eta^{-1} \cap Q \cap P\right) \backslash(\eta \Gamma \cap P)$, then $\bigsqcup_{\eta \in Y} \zeta^{-1} R_{\eta}$ gives $\bigsqcup_{\xi \in X}\left(\xi \Gamma \xi^{-1} \cap Q\right) \backslash(\xi \Gamma \cap Q P)$, where $\zeta$ is taken for each $\eta$ so that $\zeta \in Z_{\xi}$ and $\eta \in \zeta \xi \Gamma$.
(3) $\pi_{t}$ gives a bijection of $R_{\eta}$ onto $\left(\eta_{t} \Delta \eta_{t}^{-1} \cap P_{r}^{\omega}\right) \backslash \eta_{t} \Delta$, where $\Delta=\pi_{t}(\Gamma \cap P)$ and $\eta_{t}=\pi_{t}(\eta)$.

Proof. Clearly $Q \xi \Gamma=\bigsqcup_{\zeta \in Z_{\xi}}(P \cap Q) \zeta \xi \Gamma$ for every $\xi \in X$. Thus $G^{\psi}=\bigsqcup_{\xi \in X} Q \xi \Gamma$ $=\bigsqcup_{\xi, \zeta}(P \cap Q) \zeta \xi \Gamma$, and hence $P=\bigsqcup_{\xi, \zeta}(P \cap Q)(\zeta \xi \Gamma \cap P)=\bigsqcup_{\eta \in Y}(P \cap Q) \eta(\Gamma \cap P)$. Applying $\pi_{t}$ to this equality, we obtain $G^{\omega}=\bigcup_{\eta \in Y} P_{r}^{\omega} \pi_{t}(\eta(\Gamma \cap P))$ by (27.6b). To see that this union is disjoint, suppose that $\pi_{t}(\eta)=\alpha \pi_{t}\left(\eta^{\prime} \gamma\right)$ with $\alpha \in P_{r}^{\omega}, \gamma \in$ $\Gamma \cap P$ and $\eta, \eta^{\prime} \in Y$. By (27.6b) we can put $\alpha=\pi_{t}(\beta)$ with $\beta \in P \cap Q$. Then $\pi_{t}\left(\beta \eta^{\prime} \gamma \eta^{-1}\right)=1$. Since $\beta \eta^{\prime} \gamma \eta^{-1} \in P$, (27.6a) shows that $\beta \eta^{\prime} \gamma \eta^{-1} \in P \cap Q$, and hence $\eta^{\prime} \in(P \cap Q) \eta(\Gamma \cap P)$. Thus $\eta=\eta^{\prime}$, which proves the expected disjointness. Since $P_{r}^{\omega} \backslash G^{\omega} / \Gamma^{\prime}$ is finite for any congruence subgroup $\Gamma^{\prime}$ of $G^{\omega}, Y$ must be finite. Now $Q=Q^{-1}=\bigsqcup_{\zeta \in Z_{\xi}}\left(\xi \Gamma \xi^{-1} \cap Q\right) \zeta^{-1}(P \cap Q)$, and hence $Q P=\bigsqcup_{\zeta \in Z_{\xi}}\left(\xi \Gamma \xi^{-1} \cap\right.$ $Q) \zeta^{-1} P$. Therefore $\left(\xi \Gamma \xi^{-1} \cap Q\right) \backslash(\xi \Gamma \cap Q P)$ is represented by the disjoint union of $\left(\xi \Gamma \xi^{-1} \cap Q\right) \backslash\left[\xi \Gamma \cap\left(\xi \Gamma \xi^{-1} \cap Q\right) \zeta^{-1} P\right]$, which is represented by $\left(\xi \Gamma \xi^{-1} \cap Q \cap\right.$ $\left.\zeta^{-1} P \zeta\right) \backslash\left(\xi \Gamma \cap \zeta^{-1} P\right)$, ad can easily be verified. With $\eta \in \zeta \xi \Gamma \cap P$, this is clearly represented by $\zeta^{-1} R_{\eta}$. Finally we have $\pi_{t}(\eta \Gamma \cap P)=\eta_{t} \Delta$, and our last assertion can easily be verified by means of (27.6a, b).
27.8. Lemma. The notation being as in Lemma 27.7, suppose that $\lambda_{t}(\gamma)^{k}=1$ for every $\gamma \in \Gamma \cap P$; let $p_{\xi} \in \mathcal{S}_{k}^{\varphi}\left(\xi \Gamma \xi^{-1}, Q\right)$ with $\varphi=\varphi_{r}$ for each $\xi \in X$. Then

$$
\begin{align*}
& \Phi^{l-t}\left\{\delta(\mathfrak{z})^{m / 2-s \mathbf{a}} \sum_{\xi \in X} E_{k}^{\psi, \varphi}\left(\mathfrak{z}, s ; p_{\xi}, \xi \Gamma \xi^{-1}\right) \|_{k} \xi\right\}  \tag{27.7}\\
& \quad=\delta(w)^{m / 2-s \mathbf{a}} \sum_{\eta \in Y}\left|c_{\eta}\right|^{m-2 s \mathbf{a}} E_{k}^{\omega, \varphi}\left(w, s ; q_{\eta}, \eta_{t} \Delta \eta_{t}^{-1}\right) \|_{k} \eta_{t}
\end{align*}
$$

at least for $\operatorname{Re}(s)>(n+r+1) / 2$ in Case $S P$ and $\operatorname{Re}(s)>l+r+\operatorname{dim}(Z)$ in Cases $U T$ and UB. Here $c_{\eta}=\lambda_{r}\left(\zeta^{-1}\right) \lambda_{t}(\eta)$ and $q_{\eta}=c_{\eta}^{-k} p_{\xi} \|_{k} \pi_{r}\left(\zeta^{-1}\right)$ with $\zeta \in Z_{\xi}$ such that $\eta \in \zeta \xi \Gamma \cap P$.

Remark. If $\Gamma$ is sufficiently small, we have $0 \ll \lambda_{t}(\gamma) \in \mathfrak{g}^{\times}$for every $\gamma \in \Gamma \cap P$ (even if $K \neq F$ ). Then $\lambda_{t}(\gamma)^{k}=\lambda_{t}(\gamma)^{\mu \mathbf{a}}=1$. Thus the condition that $\lambda_{t}(\gamma)^{k}=1$ for every $\gamma \in \Gamma \cap P$ is satisfied for a sufficiently small $\Gamma$.

Proof. Let us first show (27.7) by formally applying $\Phi^{l-t}$ termwise. Each term is of the form $\Phi^{l-t}\left\{\delta(\mathfrak{z})^{m / 2-s \mathbf{a}}\left[\delta\left(\mathfrak{z}, s ; p_{\xi}\right) \| \alpha\right]\right\}$ with $\alpha \in\left(\xi \Gamma \xi^{-1} \cap Q\right) \backslash \xi \Gamma$. By Lemma 27.6 this is nonzero only when $\alpha \in Q P$. Thus putting $T_{\xi}=\left(\xi \Gamma \xi^{-1} \cap\right.$ $Q) \backslash(\xi \Gamma \cap Q P)$, we see that the left-hand side of (27.7) equals (formally)

$$
\begin{equation*}
\sum_{\xi \in X} \sum_{\alpha \in T_{\xi}} \Phi^{l-t}\left\{\delta(\mathfrak{z})^{m / 2-s \mathbf{a}}\left[\delta\left(\mathfrak{z}, s ; p_{\xi}\right) \| \alpha\right]\right\} . \tag{27.8}
\end{equation*}
$$

By Lemma 27.7 (2), $\bigsqcup_{\xi \in X} T_{\xi}$ can be replaced by $\bigsqcup_{\eta \in Y} \zeta^{-1} R_{\eta}$. Let $\gamma \in R_{\eta}$ and $\varepsilon=\gamma \eta^{-1}$. Then $\eta^{-1} \varepsilon \eta \in \Gamma \cap P$, and hence $\left|\lambda_{t}(\varepsilon)\right|^{\mathbf{a}}=1$ and $\lambda_{t}(\varepsilon)^{k}=1$; thus $\lambda_{t}(\gamma)^{k}=\lambda_{t}(\eta)^{k}$ and $\left|\lambda_{t}(\gamma)\right|^{\mathbf{a}}=\left|\lambda_{t}(\eta)\right|^{\mathbf{a}}$. Therefore, taking $\left(\zeta^{-1}, \gamma\right)$ as $(\beta, \gamma)$ in Lemma 27.6, we see that (27.8) equals

$$
\begin{equation*}
\delta(w)^{m / 2-s \mathbf{a}} \sum_{\eta \in Y}\left|c_{\eta}\right|^{m-2 s \mathbf{a}} \sum_{\gamma \in R_{\eta}} \delta\left(w, s ; q_{\eta}\right) \| \pi_{t}(\gamma) \tag{27.9}
\end{equation*}
$$

with $c_{\eta}$ and $q_{\eta}$ as stated in our lemma. Since $\lambda_{r}(\beta)=\lambda_{t}(\beta) \lambda_{r}\left(\pi_{t}(\beta)\right)$ for every $\beta \in P \cap Q$, we easily see that $q_{\eta} \in \mathcal{S}_{k}^{\varphi}\left(\eta_{t} \Delta \eta_{t}^{-1}, P_{r}^{\omega}\right)$. By Lemma $27.7(3), \pi_{t}\left(R_{\eta}\right)$ gives $\left(\eta_{t} \Delta \eta_{t}^{-1} \cap P_{r}^{\omega}\right) \backslash \eta_{t} \Delta$, and hence the last sum over $\gamma$ in (27.9) is $E^{\omega, \varphi}\left(w, s ; q_{\eta}\right.$, $\left.\eta_{t} \Delta \eta_{t}^{-1}\right) \| \eta_{t}$, which proves (27.7) at least in the formal sense. Now the condition $\operatorname{Re}(s)>l+r+\operatorname{dim}(Z)$ in Case UB guarantees the local uniform convergence of the series of (26.28) on $\mathfrak{Z}^{\psi}$, as well as its uniform convergence in a suitable domain $X$, as proven in [S97, Proposition A3.7]; see [S97, (A3.5.3)] for the explicit description of $X$. Now viewing the sum in question as an integral over a discrete set, we can apply the Lebesgue convergence theorem to justify our formal calculation. The necessary condition for the theorem is given by [S97, Lemma A3.6, (A3.6.4), and line 11 from the bottom on page 226]. Cases SP and UT can be handled in a similar way.
27.9. Lemma. (1) Let $\alpha \in G^{\psi}$ and $f \in \mathcal{S}_{k}^{\varphi}\left(\Gamma, P_{r}^{\psi}\right), 0 \leq r<l$; suppose $\operatorname{Re}(s)>$ $(n+r+1) / 2$ in Case $S P$ and $\operatorname{Re}(s)>l+r+\operatorname{dim}(Z)$ in Cases UT and UB. Then

$$
\Phi^{l-r}\left\{\delta(\mathfrak{z})^{m / 2-s \mathbf{a}}\left[E_{k}^{\psi, \varphi}(\mathfrak{z}, s ; f, \Gamma) \|_{k} \alpha\right]\right\}=\left\{\begin{array}{lll}
0 & \text { if } & \alpha \notin \Gamma P_{r}^{\psi} \\
\delta(w)^{m / 2-s \mathbf{a}} f & \text { if } \alpha=1
\end{array}\right.
$$

(2) If $\mathfrak{F}(\mathfrak{z})=\delta(\mathfrak{z})^{m / 2-s \mathbf{a}} \sum_{\xi \in X} E_{k}^{\psi, \varphi}\left(\mathfrak{z}, s ; p_{\xi}, \xi \Gamma \xi^{-1}\right) \|_{k} \xi$ with $X$ and $p_{\xi}$ of Lemma 27.8, then $\Phi^{l-r}\left(\mathfrak{F} \|_{k} \eta^{-1}\right)=\delta(w)^{m / 2-s \mathbf{a}} p_{\eta}$ for every $\eta \in X$ and $s$ as in (1).

Proof. Though these are essentially special cases of Lemma 27.8, it is easier to derive them directly from Lemma 27.6. In fact, the quantity of (1) is the sum of $\Phi^{l-r}\left\{\delta^{m / 2-s \mathbf{s}}\left[\delta_{f}^{\psi} \| \gamma \alpha\right]\right\}$ for $\gamma \in\left(\Gamma \cap P_{r}^{\psi}\right) \backslash \Gamma$. By Lemma 27.6 the nonvanishshing can occur only if $\gamma \alpha \in P_{r}^{\psi}$, that is, only if $\alpha \in \Gamma P_{r}^{\psi}$. This proves (1) for $\alpha \notin \Gamma P_{r}^{\psi}$. If $\alpha=1$, nonvanishshing can appear only from $\left(\Gamma \cap P_{r}^{\psi}\right) \backslash\left(\Gamma \cap P_{r}^{\psi}\right)$. which is represented by 1 . Therefore, taking $\beta=\gamma=1$ in Lemma 27.6, we obtain (1). This formal proof can be justified for the same reason as in the proof of Lemma 27.8. Assertion (2) follows immediately from (1), since $\xi \eta^{-1} \in \xi \Gamma \xi^{-1} P_{r}^{\psi}$ only if $\xi=\eta$.
27.10. Lemma. Let $0 \leq r<l$. Put $N=n+r$ in Cases $S P$ and UT and $N=l+r+\operatorname{dim}(Z)$ in Case UB; put also $\lambda(a)=(a+1) / 2$ in Case SP and $\lambda(a)=a$ in Cases UT and UB. Suppose that $\mu>\lambda(N)$ if $F \neq \mathbf{Q}$ and $\mu>\lambda(N)+1$ if $F=\mathbf{Q}$.

Suppose also that $\mu \in \Lambda(r, k)$ if $r>0$ in Cases SP and UT, and $\mu>4 r+2 \operatorname{dim}(Z)$ in Case UB, where $\Lambda(*, *)$ is defined by (23.30). Then equality (27.7) and the equalities of Lemma 27.9 are valid for $s=\mu / 2$.

Proof. Our task is to derive $\Phi^{l-t}\{\mathcal{A}(\mathfrak{z}, \mu / 2)\}=\mathcal{A}^{\prime}(w, \mu / 2)$ from the type of relation $\Phi^{l-t}\{\mathcal{A}(\mathfrak{z}, s)\}=\mathcal{A}^{\prime}(w, s)$ established in Lemmas 27.8 and 27.9. Clearly it is sufficient to prove the case $l-t=1$.

We first consider the assertion of our lemma concerning (27.7) for $t=n-1$ and $r=0$ in Cases SP and SU. This means that the desired conclusion is valid for $\mathcal{A}(z, s)=\delta(z)^{m / 2-s \mathbf{a}}\left[E_{k}^{n, 0}(z, s ; 1, \Gamma) \| \alpha\right], \alpha \in G^{n}$. The proof given in [S95a, pp.579-580] is quite technical and requires rather involved preliminaries [S95a, pp.549-551], and so we refer the reader to those pages for the detailed proof. The fact was proven only in Case SP, but the proof is applicable to Case SU with obvious modifications. (We can easily give the analogues of Lemmas 2.3, 2.4, and 2.6 of [S95a] in Case SU. In Case UT we may assume that $\Gamma$ is contained in $\operatorname{SU}\left(\eta_{n}\right)$, as explained in the proof of Lemma 17.13; the condition $\operatorname{Re}(2 s)>n+1$ should be changed for $\operatorname{Re}(s)>n$. In the proof we need Proposition 6.16, as well as the holomorphy of $E_{\mu \mathbf{a}}^{n, 0}(z, \mu / 2 ; \mu \mathbf{a}, \Gamma)$ and $E_{\mu \mathbf{a}}^{n-1,0}(z, \mu / 2 ; \mu \mathbf{a}, \Gamma)$ which is guaranteed by Theorem 17.7 (i). Also we have to assume that $\mu \neq \lambda(n)$ (see [S95a, p580, line $11]$ ), which is why we have to assume $\mu>\lambda(n)$ instead of $\mu \geq \lambda(n)$ even if $F \neq \mathbf{Q}$.)

Let us now consider the action of $\Phi$ on $E_{k}^{\psi, \varphi}$ in the general case. Let $\mathcal{A}(\mathfrak{z}, s)$ denote the function inside the brackets of (27.7). We first take Case UB. By [S97, Proposition 20.10] $\mathcal{A}(\mathfrak{z}, s)$ is a finite linear combination of functions of the form $\delta(\mathfrak{z})^{m / 2-s \mathbf{a}}\left[E_{1}\left(\mathfrak{z}, s ; \mathbf{f}, \chi, D^{\psi}\right) \|_{k} \alpha\right]$ with $\alpha \in G^{\psi}$. We may assume that each $\mathbf{f}$ is an eigenform of Hecke operators, since $\mathcal{S}_{k}^{\varphi}\left(D^{\varphi}\right)$ with $D^{\varphi}$ of (26.32) is spanned by such eigenforms (see [S97, Proposition 20.4 (1)]). Taking $p=1$ in (26.42) and applying $\| \alpha$, we obtain, employing [S97, (22.3.3)],

$$
\begin{aligned}
& \mu c_{m}(\mathbf{s}) C^{\prime}(s) \mathcal{Z}(s, \mathbf{f}, \chi) E_{1}\left(\mathfrak{z}, s ; \mathbf{f}, \chi, D^{\psi}\right) \|_{k} \alpha \\
& \quad=\Lambda_{\mathfrak{c}}^{n}(s, \chi) \sum_{a \in \mathcal{B}} \chi_{\mathbf{h}}(\operatorname{det}(a)) \int_{\mathfrak{D}_{a}}\left(H_{1, a} \|_{k}[\alpha, 1]_{S}\right)^{\circ}(\mathfrak{z}, w ; s) f_{a}(w) \delta(w)^{m} \mathbf{d} w,
\end{aligned}
$$

where $n=\operatorname{dim}\left(V_{r}\right)$. Recall that $\mathcal{F}^{\circ}(\mathfrak{z}, w)=\delta\left(w, \wp_{r}(\mathfrak{z})\right)^{-m} \mathcal{F}\left(\iota_{U}(\mathfrak{z}, w)\right)$ for a function $\mathcal{F}$ on $\mathcal{H}^{N}$, where $N=l+r+\operatorname{dim}(Z)$, and $H_{1, a}(z)$ for $z \in \mathcal{H}^{N}$ corresponds to $E_{q_{0}}(z)$ with some $q_{0} \in G_{\mathrm{h}}^{N}$ in the sense of (17.23a). Put $\mathfrak{E}_{a}=H_{1, a} \|_{k}[\alpha, 1]_{S}$ and $\mathfrak{F}_{a}(z, s)=\delta(z)^{m / 2-s \mathbf{a}} \mathfrak{E}_{a}(z, s)$ for $z \in \mathcal{H}^{N}$. By Lemma 17.13, we can put

$$
\begin{equation*}
\mathfrak{E}_{a}(z, s)=\sum_{i \in I} b_{i} c_{i}^{s} E\left(z, s ; k, \Gamma_{i}\right) \|_{k} \alpha_{i} \tag{}
\end{equation*}
$$

with a finite set $I$ and $b_{i}, c_{i}, \Gamma_{i}, \alpha_{i}$ as described there. The map $\iota_{U}: \mathfrak{Z}^{\psi} \times \mathfrak{Z}^{\varphi} \rightarrow \mathcal{H}^{N}$ is defined by [S97, (22.2.1)]. Put $\omega=\varphi_{l-1}$ nd define $\iota_{U_{1}}: \mathfrak{Z}^{\omega} \times \mathfrak{Z}^{\varphi} \rightarrow \mathcal{H}^{N-1}$ in the same way, where $U_{1}$ denotes the element of $U\left(\eta_{N-1}\right)_{\mathbf{a}}$ corresponding to $U$. Then we can easily verify that $\iota_{U}\left(\operatorname{diag}\left[\rho \mathbf{i}_{1}, \mathfrak{z}_{1}\right], w\right)=\operatorname{diag}\left[\rho \mathbf{i}_{1}, \iota_{U_{1}}\left(\mathfrak{z}_{1}, w\right)\right]$ for $\mathfrak{z}_{1} \in \mathfrak{Z}^{\omega}$ and $w \in \mathfrak{Z}^{\varphi}$. (Relevant formulas are [S97, (6.10.2), (22.1.2), (22.1.3), (22.1.4), (22.1.6), and (22.2.8)]. From these we can easily derive that $U=1_{2} \times U_{1}$ in the sense of (23.5).) Also $\wp_{r}^{\psi}\left(\operatorname{diag}\left[\boldsymbol{i}_{1}, \mathfrak{z}_{1}\right]\right)=\wp_{r}^{\omega}\left(\mathfrak{z}_{1}\right)$.

Put $\Phi\left(\mathfrak{F}_{a}(z, s)\right)=\mathfrak{F}_{a}^{\prime}\left(z_{1}, s\right)$, where $z_{1} \in \mathcal{H}^{N-1}$. Then $\mathfrak{F}_{a}^{\prime}\left(z_{1}, s\right)=\delta\left(z_{1}\right)^{m / 2-s \mathbf{a}}$ $\cdot \mathfrak{E}_{a}^{\prime}\left(z_{1}, s\right)$ for $\operatorname{Re}(s)>\lambda(N)$ with a function $\mathfrak{E}_{a}^{\prime}$ of type $\left(^{*}\right)$ defined for $z_{1} \in \mathcal{H}^{N-1}$. This is a special case of Lemma 27.8. Moreover, $\Phi\left(\mathfrak{F}_{a}(z, \mu / 2)\right)=\mathfrak{F}_{a}^{\prime}\left(z_{1}, \mu / 2\right)$ for
$\mu>\lambda(N)$ if $F \neq \mathbf{Q}$ and $\mu>\lambda(N)+1$ if $F=\mathbf{Q}$, as we said in the second paragraph of this proof. Now

$$
\mathfrak{F}_{a}^{\circ}\left(\operatorname{diag}\left[\rho \mathbf{i}_{1}, \mathfrak{z}_{1}\right], w ; s\right)=\delta\left(w, \wp_{r}\left(\mathfrak{z}_{1}\right)\right)^{-m} \mathfrak{F}_{a}\left(\operatorname{diag}\left[\rho \mathbf{i}_{1}, \iota_{U_{1}}\left(\mathfrak{z}_{1}, w\right)\right], s\right)
$$

so that $\left(\Phi\left(\mathfrak{F}_{a}\right)^{\circ}\right)\left(\mathfrak{z}_{1}, w ; s\right)=\left(\mathfrak{F}_{a}^{\prime}\right)^{\circ}\left(\mathfrak{z}_{1}, w ; s\right)$, which holds for $\operatorname{Re}(s)>N$ and also for $s=\mu / 2$ as above. By $[\operatorname{S97},(22.2 .7)]$ we have

$$
\delta\left(\iota_{U}(\mathfrak{z}, w)\right)=h\left|\delta\left(w, \wp_{r}(\mathfrak{z})\right)\right|^{-2} \delta(\mathfrak{z}) \delta(w)
$$

with a positive costant $h$, and the formula holds for $\left(\mathfrak{z}_{1}, U_{1}\right)$ in place of $(\mathfrak{z}, U)$ with the same $h$. Put $\mathfrak{A}(\mathfrak{z}, s)=\delta(\mathfrak{z})^{m / 2-s \mathbf{a}} E_{1}\left(\mathfrak{z}, s ; \mathbf{f}, \chi, D^{\psi}\right) \|_{k} \alpha$. Then
(**) $\quad \mathcal{D}(s) \mathfrak{A}(\mathfrak{z}, s)$

$$
=h^{s \mathbf{a}-m / 2} \sum_{a \in \mathcal{B}} e_{a} \int_{\mathfrak{D}_{a}}\left(\mathfrak{F}_{a}\right)^{\circ}(\mathfrak{z}, w ; s)\left|\delta\left(w, \wp_{r}(\mathfrak{z})\right)\right|^{m-2 s \mathbf{a}} f_{a}(w) \delta(w)^{m / 2+s \mathbf{a}} \mathbf{d} w
$$

with $e_{a} \in \mathbf{C}$ and $\mathcal{D}(s)=c_{m}(\mathbf{s}) C^{\prime}(s) \mathcal{Z}(s, \mathbf{f}, \chi) \Lambda_{\mathrm{c}}^{n}(s, \chi)^{-1}$.
To see the behavior of the integrals under $\Phi$, take $\mathfrak{z}=\operatorname{diag}\left[\rho \mathbf{i}_{1}, \mathfrak{z}_{1}\right]$ with $\rho>1$ and $\mathfrak{z}_{1}$ in a compact subset of $\mathfrak{Z}^{\omega}$. Since $\iota(\mathfrak{z}, w)=\operatorname{diag}\left[\rho \mathbf{i}_{1}, \iota\left(\mathfrak{z}_{1}, w\right)\right]$, we see that $\iota(\mathfrak{z}, w)$ belongs to the subset $\mathfrak{X}$ of $\mathcal{H}^{N}$ described in [S97, p.189, line 7 from the bottom] if $w$ belongs to the Siegel set $\mathfrak{S}^{\prime}$ employed there. Therefore the argument of [S97, §22.12] proves the absolute and uniform convergence of the integrals of ${ }^{(* *)}$ for such $\mathfrak{z}$ and $s$ belonging to a compact subset of C. (In [S97, §22.12], $\mathfrak{z}$ was taken in a compact set, but the argument there is valid for any $\mathfrak{z}$ such that $\iota\left(\mathfrak{z}, \mathfrak{S}^{\prime}\right) \subset \mathfrak{X}$.) Thus we can apply $\Phi$ to $\left({ }^{(* *)}\right.$ and obtain

$$
\begin{aligned}
\left(^{* * *}\right) & \mathcal{D}(s)(\Phi \mathfrak{A})\left(\mathfrak{z}_{1}, s\right) \\
= & h^{s \mathbf{a}-m / 2} \sum_{a \in \mathcal{B}} e_{a} \int_{\mathfrak{D}_{a}}\left(\mathfrak{F}_{a}^{\prime}\right)^{\circ}\left(\mathfrak{z}_{1}, w ; s\right)\left|\delta\left(w, \wp_{r}\left(\mathfrak{z}_{1}\right)\right)\right|^{m-2 s \mathbf{a}} f_{a}(w) \delta(w)^{m / 2+s \mathbf{a}} \mathbf{d} w \\
= & \delta\left(\mathfrak{z}_{1}\right)^{m / 2-s \mathbf{a}} \sum_{a \in \mathcal{B}} e_{a} \int_{\mathfrak{D}_{a}}\left(\mathfrak{E}_{a}^{\prime}\right)^{\circ}\left(\mathfrak{z}_{1}, w ; s\right) f_{a}(w) \delta(w)^{m} \mathbf{d} w
\end{aligned}
$$

The last integrals define meromorphic functions of $s$ on the whole $\mathbf{C}$; moreover they are meaningful at $s=\mu / 2$. Now $\mathcal{D}$ is finite and nonzero at $s=\mu / 2$, as already seen in the proof of Theorerm 23.11 given in $\S 25.7$ in Cases SP and UT and in the proof of Theorem 26.13 in Case UB. Thus $\Phi(\mathfrak{A}(*, \mu / 2))$ is meaningful and equals $\mathcal{D}(\mu / 2)^{-1}$ times the last line of $\left({ }^{* * *}\right)$ at $s=\mu / 2$.

Returning to the question at the beginning, let $\mathcal{A}^{\prime}\left(\mathfrak{z}_{1}, s\right)$ denote the right-hand side of $(27.7)$ with $l-t=1$. Then $(\Phi \mathcal{A})\left(\mathfrak{z}_{1}, s\right)=\mathcal{A}^{\prime}\left(\mathfrak{z}_{1}, s\right)$ for $\operatorname{Re}(s)>\lambda(N)$. Since $\mathcal{A}$ is a finite linear combination of functions of type $\mathfrak{A}$, we see that $\mathcal{A}^{\prime}\left(\mathfrak{z}_{1}, s\right)$ is a inear combination of some functions that can be given by the last line of ( ${ }^{* * *}$ ) times $\mathcal{D}(s)^{-1}$. Therefore $\Phi(\mathcal{A}(*, \mu / 2))$ is meaningful and must coincide with $\mathcal{A}^{\prime}\left(\mathfrak{z}_{1}, \mu / 2\right)$. This proves the assertion of our lemma concerning (27.7) in Case UB for $l-t=1$, which proves, by induction, the general case. Cases SP and UT are similar. As for the assertion concerning the equalities of Lemma 27.9, the left-hand sides are special cases of that of (27.7); the only point is that the right-hand sides can be given explicitly as stated in that lemma. Therefore the above proof is applicable to them. This completes the proof.
27.11. With $k, m$, and $\mu$ as in $\S 27.1$, for $0 \leq r \leq l$ we put, in Case UB,

$$
\begin{equation*}
E_{k}^{\psi, r}(z ; f, \Gamma)=E_{k}^{\psi, \varphi_{r}}(z, \mu / 2 ; f, \Gamma) \quad\left(z \in \mathfrak{Z}^{\psi}\right) \tag{27.10}
\end{equation*}
$$

for $f \in \mathcal{S}_{k}^{\varphi_{r}}\left(\Gamma, P_{r}^{\psi}\right)$, whenever the right-hand side is finite. (Hereafter we use $z$ instead of $\mathfrak{z}$ for the variable on $\mathfrak{Z}^{\psi}$.) In general this may not be holomorphic in $z$. We denote by $\tilde{\mathcal{E}}_{k}^{\psi, r}$ the vector space spanned over $\mathbf{C}$ by $E_{k}^{\psi, r}(z ; f, \Gamma) \|_{k} \alpha$ for all $\alpha \in G^{\psi}$, all congruence subgroups $\Gamma$ of $G^{\psi}$, and all such $f$; we then put $\mathcal{E}_{k}^{\psi, r}=\tilde{\mathcal{E}}_{k}^{\psi, r} \cap \mathcal{M}_{k}^{\psi}$. Clearly these spaces are stable under the operator $\|_{k} \xi$ for every $\xi \in G^{\psi}$. In Cases SP and UT we take $E_{k}^{n, r}(z, s ; f, \Gamma)$ as $E_{k}^{\psi, r}(z, s ; f, \Gamma)$. Therefore the symbols $\tilde{\mathcal{E}}_{k}^{n, r}$ and $\mathcal{E}_{k}^{n, r}$ are more natural in those two cases, but we state our results in all three cases by using $\tilde{\mathcal{E}}_{k}^{\psi, r}$ and $\mathcal{E}_{k}^{\psi, r}$, and give the proof mainly in Case UB, although we indicate necessary modifications in the other two cases. The reader is reminded of the remark we made in $\S 27.1$; for example, $l=n$ and $\psi=\eta_{n}$ in Cases SP and UT; the case of half-integral weight can be included.

If $r=l$, then $E_{k}^{\psi, \psi}(z, s ; f, \Gamma)=f(z)$ for every $f \in \mathcal{S}_{k}^{\psi}(\Gamma)$ by (23.14) and the remark at the end of $\S 26.9$, so that

$$
\begin{equation*}
\tilde{\mathcal{E}}_{k}^{\psi, l}=\mathcal{E}_{k}^{\psi, l}=\mathcal{S}_{k}^{\psi} . \tag{27.11}
\end{equation*}
$$

We have $k=m=\mu \mathbf{a}$ in Case SP. Now, in Cases UT and UB we can restrict $\Gamma$ to the congruence subgroups of $S U(\psi)$, by virtue of (23.13a) and (26.28a). Since $j_{\gamma}^{k}=c_{\gamma} j_{\gamma}^{\mu \mathbf{a}}$ for $\gamma \in G^{\psi}$ with $c_{\gamma} \in \overline{\mathbf{Q}}^{\times}$and $c_{\gamma}=1$ for $\gamma \in S U(\psi)$, the spaces $\tilde{\mathcal{E}}_{k}^{\psi, r}$ and $\mathcal{E}_{k}^{\psi, r}$ depend only on $\mu$. We have to be more careful, however, when we speak of the rationality of modular forms over a number field.
27.12. Lemma. Let $0 \leq r \leq l$. Impose the condition of Lemma 27.10 on $\mu$ if $r<l$. Then the function of (27.10) is meaningful and the following assertions hold:
(1) $\mathcal{E}_{k}^{\psi, r}=\widetilde{\mathcal{E}}_{k}^{\psi, r}$.
(2) $\Phi^{l-t}\left(\mathcal{E}_{k}^{\psi, r}\right) \subset \mathcal{E}_{k}^{\varphi_{t}, r}$ for $r \leq t \leq l$ and in particular $\Phi^{l-r}\left(\mathcal{E}_{k}^{\psi, r}\right)=\mathcal{S}_{k}^{\varphi_{r}}$.
(3) $\Phi^{s}\left(\mathcal{E}_{k}^{\psi, r}\right)=0$ if $s>l-r$.
(4) If $g \in \mathcal{E}_{k}^{\psi, r}$ and $\Phi^{l-r}\left(g \|_{k} \alpha\right)=0$ for every $\alpha \in G^{\psi}$, then $g=0$.

Proof. All the assertions are trivial if $r=l$, and so we assume $r<l$; notice that (3) for $r=l$ follows from Lemma 27.2 (2). Taking $m=\mu$ a in Theorems 17.7 (i), $23.11(\mathrm{I}), 26.13(1)$, we see that $E_{k}^{\psi, r}(z ; f, \Gamma)$ is meaningful and holomorphic. Thus we obtain (1). Assertion (2) follows from Lemmas 27.8, 27.9, and 27.10. Then $\Phi^{l-r+1}\left(\mathcal{E}_{k}^{\psi, r}\right)=\Phi\left(\Phi^{l-r}\left(\mathcal{E}_{k}^{\psi, r}\right)\right)=\Phi\left(\mathcal{S}_{k}^{\varphi_{r}}\right)=0$ by Lemma 27.2 (2), which gives (3). To prove (4), let $g=\sum_{i \in I} E_{k}^{\psi, r}\left(z ; f_{i}, \Gamma_{i}\right) \| \alpha_{i}$ with a finite set of indices $I, \alpha_{i} \in G^{\psi}$, congruence subgroups $\Gamma_{i}$, and $f_{i} \in \mathcal{S}_{k}^{\varphi_{r}}\left(\Gamma_{i}, P_{r}^{\psi}\right)$. Take a congruence subgroup $\Gamma$ so that $\Gamma \subset \bigcap_{i \in I} \alpha_{i}^{-1} \Gamma_{i} \alpha_{i}$ and $\lambda_{t}\left(\Gamma \cap \mathcal{P}_{t}^{\psi}\right)^{k}=1$ for every $t \geq r$. For each $i$ we can find a finite set $B_{i}$ such that $\Gamma_{i} \alpha_{i}=\bigsqcup_{\beta \in B_{i}} \beta \Gamma$. Then $\beta \Gamma \beta^{-1} \subset \Gamma_{i}$ for every $\beta \in B_{i}$, and $\Gamma_{i}=\bigsqcup_{\beta \in B_{i}} \beta \Gamma \beta^{-1} \beta \alpha_{i}^{-1}$. Therefore, taking $\Gamma_{i}, \beta \Gamma \beta^{-1}$, and $\beta \alpha_{i}^{-1}$ as $\Gamma, \Gamma^{\prime}$, and $\alpha$ in (23.13a) or (26.28a), we find that

$$
g=\sum_{i \in I} \sum_{\beta \in B_{i}} E_{k}^{\psi, r}\left(z ; c_{i, \beta} f_{i}, \beta \Gamma \beta^{-1}\right) \| \beta
$$

with some constants $c_{i, \beta}$. Thus we may assume at the beginning that $\Gamma_{i}=\alpha_{i} \Gamma \alpha_{i}^{-1}$ for every $i \in I$ with some $\Gamma$. Now if $\beta \in \sigma \xi \Gamma$ with $\sigma \in P_{r}^{\psi}$ and $\xi \in G^{\psi}$, then $\left(\beta \Gamma \beta^{-1} \cap P_{r}^{\psi}\right) \backslash \beta \Gamma$ can be given by $\sigma T$ with $T=\left(\xi \Gamma \xi^{-1} \cap P_{r}^{\psi}\right) \backslash \xi \Gamma$. Also, for $\tau \in T$ we have

$$
\delta_{s . f}\left\|\sigma \tau=\lambda_{r}(\sigma)^{-k}\left|\lambda_{r}(\sigma)\right|^{m-2 s \mathbf{a}} \delta\left(z, s ; f \| \pi_{r}(\sigma)\right)\right\| \tau
$$

by (23.12) or [S97, (12.3.5)], and hence

$$
E_{k}^{\psi, r}\left(z ; f, \beta \Gamma \beta^{-1}\right)\left\|\beta=\lambda_{r}(\sigma)^{-k} E_{k}^{\psi, r}\left(z ; f \| \pi_{r}(\sigma), \xi \Gamma \xi^{-1}\right)\right\| \xi
$$

Therefore, with $X=P_{r}^{\psi} \backslash G^{\psi} / \Gamma$ we can put $g=\sum_{\xi \in X} E_{k}^{\psi, r}\left(z ; h_{\xi}, \xi \Gamma \xi^{-1}\right) \| \xi$. Suppose that $\Phi^{l-r}(g \| \alpha)=0$ for every $\alpha \in G^{\psi}$. Then by Lemmas 27.9 (2) and 27.10 we have $h_{\eta}=\Phi^{l-r}\left(g \| \eta^{-1}\right)=0$ for every $\eta \in X$, so that $g=0$, which proves (4).

We are now ready to state our main theorems on the structure of the spaces of holomorphic Eisenstein series.
27.13. Theorem. Suppose $n>1$ in Cases $S P$ and UT, and $2 l+\operatorname{dim}(Z)>2$ in Case UB; put $\mathcal{E}_{k}^{\psi, r}(\Gamma)=\mathcal{E}_{k}^{\psi, r} \cap \mathcal{M}_{k}^{\psi}(\Gamma)$. Then we have

$$
\mathcal{M}_{k}^{\psi}=\bigoplus_{r=0}^{l} \mathcal{E}_{k}^{\psi, r} \quad \text { and } \quad \mathcal{M}_{k}^{\psi}(\Gamma)=\bigoplus_{r=0}^{l} \mathcal{E}_{k}^{\psi, r}(\Gamma)
$$

for every congruence subgroup $\Gamma$ of $G^{\psi}$ provided the following condition is satisfied:
(27.12) Case SP: $\mu \geq 3 n / 2$ if $n>2 ; \mu>3$ if $n=2$ and $F=\mathbf{Q} ; \mu>2$ if $n=2$ and $F \neq \mathbf{Q}$;
Case UT: $\mu>4$ if $F=\mathbf{Q}$ and $n=2 ; \mu \geq 3 n-2$ otherwise;
Case UB: $\mu>3$ if $F=\mathbf{Q}$ and $l=\operatorname{dim}(Z)=1 ; \mu \geq 2 \operatorname{dim}(W)-3$ otherwise.

Proof. We need the condition on $\mu$ of Lemma 27.10 for every $r<l$, which is why (27.12) is required. Suppose that $\sum_{r=0}^{l} p_{r}=0$ with $p_{r} \in \mathcal{E}_{k}^{\psi, r}$. By Lemma 27.12 (3) we have $\Phi^{l}\left(p_{0} \| \alpha\right)=-\sum_{r=1}^{l} \Phi^{l}\left(p_{r} \| \alpha\right)=0$ for every $\alpha \in G^{\psi}$, and hence $p_{0}=0$ by Lemma 27.12 (4). Similarly we find that $\Phi^{l-1}\left(p_{1} \| \alpha\right)=0$ for every $\alpha \in G^{\psi}$, which means that $p_{1}=0$ for the same reason. Repeating this process, we obtain $p_{r}=0$ for every $r$, which proves that the $\mathcal{E}_{k}^{\psi, r}$ for $0 \leq r \leq l$ form a direct sum. Now given $g \in \mathcal{M}_{k}^{\psi}$, take $\Gamma$ so that $g \in \mathcal{M}_{k}^{\psi}(\Gamma)$ and $\lambda_{r}(\Gamma \cap$ $\left.P_{r}^{\psi}\right)^{k}=1$ for every $r$; take also $X_{r}=P_{r}^{\psi} \backslash G^{\psi} / \Gamma$. Put $h_{\xi}=\Phi^{l}\left(g \| \xi^{-1}\right)$ for each $\xi \in X_{0}$ and $f_{0}=\sum_{\xi \in X_{0}} E_{k}^{\psi, 0}\left(z ; h_{\xi}, \xi \Gamma \xi^{-1}\right) \| \xi$. Then by Lemma 27.9 (2) and Lemma 27.10, $\Phi^{l}\left(\left(g-f_{0}\right) \| \xi^{-1}\right)=0$ for every $\xi \in X_{0}$. By Lemma 27.3 we have $\Phi^{l-1}\left(\left(g-f_{0}\right) \| \alpha\right) \in \mathcal{S}_{k}^{\varphi_{1}}$ for every $\alpha \in G^{\psi}$. Put $p_{\eta}=\Phi^{l-1}\left(\left(g-f_{0}\right) \| \eta^{-1}\right)$ for $\eta \in X_{1}$, and $f_{1}=\sum_{\eta \in X_{1}} E_{k}^{\psi, 1}\left(z ; p_{\eta}, \eta \Gamma \eta^{-1}\right) \| \eta$. Then $\Phi^{l-1}\left(\left(g-f_{0}-f_{1}\right) \| \eta^{-1}\right)=0$ for every $\eta \in X_{1}$, and hence $\Phi^{l-2}\left(\left(g-f_{0}-f_{1}\right) \| \alpha\right) \in \mathcal{S}_{k}^{\varphi_{2}}$ for every $\alpha \in G^{\psi}$ by Lemma 27.3. Continuing in this fashion, we find some elements $f_{r} \in \mathcal{E}_{k}^{\psi, r}$ for $r \leq l-1$ so that if we put $f_{l}=g-\sum_{r=0}^{l-1} f_{r}$, then $\Phi\left(f_{l} \| \alpha\right)=0$ for every $\alpha \in G^{\psi}$, which means that $f_{l} \in \mathcal{S}_{k}^{\psi}=\mathcal{E}_{k}^{\psi, l}$. This proves the first equality. If $g \in \mathcal{M}_{k}^{\psi}(\Gamma)$, then $\sum_{r=0}^{l} f_{r} \| \gamma=\sum_{r=0}^{l} f_{r}$ for every $\gamma \in \Gamma$, so that $f_{r} \| \gamma=f_{r}$, that is, $f_{r} \in \mathcal{M}_{k}^{\psi}(\Gamma) \cap \mathcal{E}_{k}^{\psi, r}$, which completes the proof.
27.14. Theorem. Let $n$ and $\mu$ be as in Theorem 27.13; put $\mathcal{E}_{k}^{\psi}=\sum_{r=0}^{l-1} \mathcal{E}_{k}^{\psi, r}$ and $\mathcal{E}_{k}^{\psi}(\Gamma)=\mathcal{E}_{k}^{\psi} \cap \mathcal{M}_{k}^{\psi}(\Gamma)$. Then

$$
\begin{aligned}
& \mathcal{M}_{k}^{\psi}=\mathcal{S}_{k}^{\psi} \oplus \mathcal{E}_{k}^{\psi}, \quad \mathcal{E}_{k}^{\psi}(\Gamma)=\bigoplus_{r=0}^{l-1} \mathcal{E}_{k}^{\psi, r}(\Gamma), \\
& \mathcal{E}_{k}^{\psi}=\left\{f \in \mathcal{M}_{k}^{\psi} \mid\langle f, g\rangle=0 \text { for every } g \in \mathcal{S}_{k}^{\psi}\right\}, \\
& \mathcal{E}_{k}^{\psi, r}=\left\{f \in \mathcal{E}_{k}^{\psi} \mid \Phi\left(f \|_{k} \alpha\right) \in \mathcal{E}_{k}^{\varphi_{l-1}, r} \text { for every } \alpha \in G^{\psi}\right\} \quad \text { if } \quad r<l .
\end{aligned}
$$

Moreover, if $\mathfrak{q}: \mathcal{M}_{k}^{\psi} \rightarrow \mathcal{S}_{k}^{\psi}$ denotes the projection map determined by the decomposition $\mathcal{M}_{k}^{\psi}=\mathcal{S}_{k}^{\psi} \oplus \mathcal{E}_{k}^{\psi}$, then $\mathfrak{q}(f)^{\sigma}=\mathfrak{q}\left(f^{\sigma}\right)$ for every $\sigma \in \operatorname{Aut}(\mathbf{C} / \overline{\mathbf{Q}})$.

Proof. The first and second equalities follow immediately from Theorem 27.13. If $f \in \mathcal{E}_{k}^{\psi, r}, r<l$, and $g \in \mathcal{S}_{k}^{\psi}$, then $\langle f, g\rangle=0$ by [S97, Lemma A3.8]. (Take the present $\delta_{s, f}, P_{r}^{\psi}, 1$ to be $g, P, \tau$ in that lemma.) This together with the first equality gives the third equality. Put $\omega=\varphi_{l-1}$; suppose $r<l$; then by Lemma $27.12(2), \Phi(f \| \alpha) \in \mathcal{E}_{k}^{\omega, r}$ for every $f \in \mathcal{E}_{k}^{\psi, r}$ and every $\alpha \in G^{\psi}$. Conversely suppose that $f=\sum_{s=0}^{l-1} g_{s}$ with $g_{s} \in \mathcal{E}_{k}^{\psi, s}$ and that $\Phi(f \| \alpha) \in \mathcal{E}_{k}^{\omega, r}$ for every $\alpha \in G^{\psi}$. Since $\Phi\left(g_{s} \| \alpha\right) \in \mathcal{E}_{k}^{\omega, s}$ by Lemma 27.12 (2) and the $\mathcal{E}_{k}^{\omega, s}$ for $0 \leq s<l$ form a direct sum, we obtain $\Phi\left(g_{s} \| \alpha\right)=0$ for $s \neq r$. By Lemma 27.2 (2), $g_{s} \in \mathcal{E}_{k}^{\psi, s} \cap \mathcal{S}_{k}^{\psi}=\{0\}$ for $s \neq r$, and hence $f=g_{r} \in \mathcal{E}_{k}^{\psi, r}$, which proves the third equality. The last assertion concerning $\mathfrak{q}$ will be proven at the end of the proof of Theorem 27.16.
27.15. Theorem. Let $0 \leq r<l$. Put $N=n+r$ in Cases $S P$ and UT and $N=l+r+\operatorname{dim}(Z)$ in Case UB; put also $\lambda(a)=(a+1) / 2$ in Case SP and $\lambda(a)=a$ in Cases UT and UB. Suppose that $\mu>\lambda(N)$ if $F \neq \mathbf{Q}$ and $\mu>\lambda(N)+1$ if $F=\mathbf{Q}$. Suppose also that $\mu \in \Lambda(r, k)$ if $r>0$ in Cases SP and UT, and $\mu>4 r+2 \operatorname{dim}(Z)$ in Case UB, where $\Lambda(*, *)$ is defined by (23.30). Given a congruence subgroup $\Gamma$ of $G^{\psi}$, let $\mathcal{E}_{k}^{\psi, r}(\Gamma)=\mathcal{M}_{k}(\Gamma) \cap \mathcal{E}_{k}^{\psi, r}$ and let $X=P_{r}^{\psi} \backslash G^{\psi} / \Gamma$ with a fixed $r<$ l. Then $f \mapsto\left(\Phi^{l-r}\left(f \| \xi^{-1}\right)\right)_{\xi \in X}$ gives a $\mathbf{C}$-linear isomorphism of $\mathcal{E}_{k}^{\psi, r}(\Gamma)$ onto $\prod_{\xi \in X} \mathcal{S}_{k}^{\varphi_{r}}\left(\xi \Gamma \xi^{-1}, P_{r}^{\psi}\right)$. Moreover, if $f \in \mathcal{E}_{k}^{\psi, r}(\Gamma)$ and $p_{\xi}=\Phi^{l-r}\left(f \| \xi^{-1}\right)$, then $f=\sum_{\xi \in X} E_{k}^{\psi, r}\left(z ; p_{\xi}, \xi \Gamma \xi^{-1}\right) \|_{k} \xi$.

Proof. If $f \in \mathcal{E}_{k}^{\psi, r}(\Gamma)$, then $\Phi^{l-r}\left(f \| \xi^{-1}\right) \in \mathcal{S}_{k}^{\varphi_{r}}\left(\xi \Gamma \xi^{-1}, P_{r}^{\psi}\right)$ by Lemma 27.2 (4) and Lemma 27.12 (2), and so our map is meaningful. The injectivity of the map follows from Lemma 27.3 and Lemma 27.12 (4). Now, given $p_{\xi} \in \mathcal{S}_{k}^{\varphi_{r}}\left(\xi \Gamma \xi^{-1}, P_{r}^{\psi}\right)$ for each $\xi \in X$, put $g=\sum_{\xi \in X} E_{k}^{\psi, r}\left(z ; p_{\xi}, \xi \Gamma \xi^{-1}\right) \|_{k} \xi$. Then $g \in \mathcal{E}_{k}^{\psi, r}(\Gamma)$ and $\Phi^{l-r}\left(g \| \xi^{-1}\right)=p_{\xi}$ for every $\xi \in X$ by Lemma 27.10 , which proves the surjectivity.
27.16. Theorem. Let $\mu$ and $r$ be as in Theorem 27.15. Put $\mathcal{E}_{k}^{\psi, r}(\Gamma, \overline{\mathbf{Q}})=$ $\mathcal{E}_{k}^{\psi, r} \cap \mathcal{M}_{k}^{\psi}(\Gamma, \overline{\mathbf{Q}})$ and $\mathcal{E}_{k}^{\psi}(\Gamma, \overline{\mathbf{Q}})=\mathcal{E}_{k}^{\psi} \cap \mathcal{M}_{k}^{\psi}(\Gamma, \overline{\mathbf{Q}})$. Then
(1) $E_{k}^{\psi, r}(z ; f, \Gamma)$ is $\overline{\mathbf{Q}}$-rational if $f$ is $\overline{\mathbf{Q}}$-rational.
(2) $\mathcal{E}_{k}^{\psi, r}(\Gamma)=\mathcal{E}_{k}^{\psi, r}(\Gamma, \overline{\mathbf{Q}}) \otimes_{\overline{\mathbf{Q}}} \mathbf{C}$.
(3) $\mathcal{E}_{k}^{\psi}(\Gamma)=\mathcal{E}_{k}^{\psi}(\Gamma, \overline{\mathbf{Q}}) \otimes_{\overline{\mathbf{Q}}} \mathbf{C}$ if (27.12) is satisfied.

Proof. Assertion (1) will be proven in $\S \S 28.12$ and 29.9. Now every element $g$ of $\mathcal{E}_{k}^{\psi, r}(\Gamma)$ can be written $g=\sum_{\xi \in X} E_{k}^{\psi, r}\left(z ; q_{\xi}, \xi \Gamma \xi^{-1}\right) \| \xi$ with $q_{\xi} \in$ $\mathcal{S}_{k}^{\varphi_{r}}\left(\xi \Gamma \xi^{-1}, P_{r}^{\psi}\right)$ by Theorem 27.15. By Lemma 23.13 and (26.25), $q_{\xi}$ is a C-linear combination of elements of $\mathcal{S}_{k}^{\varphi_{r}}\left(\xi \Gamma \xi^{-1}, P_{r}^{\psi}, \overline{\mathbf{Q}}\right)$. Therefore (2) follows from (1). Assertion (3) is immediate from (2) and the second equality of Theorem 27.14.

To prove the last assertion of Theorem 27.14, take $\sigma \in \operatorname{Aut}(\mathbf{C} / \overline{\mathbf{Q}})$. Then $\left(\mathcal{E}_{k}^{\psi}\right)^{\sigma}=$ $\mathcal{E}_{k}^{\psi}$ by (3); also $\left(\mathcal{S}_{k}^{\psi}\right)^{\sigma}=\mathcal{S}_{k}^{\psi}$ by Theorem 10.8 (1) and Proposition 26.8 (2). The desired fact now follows immediately from these equalities.
27.17. Remark. (A) We excluded the case $n=1$ in Theorems 27.13 and 27.14. If $n=1$ in Case UT, we have $S U\left(\eta_{1}\right)=S L_{2}(F)$. Thus the case $n=1$ is about the Eisenstein series for $S L_{2}(F)$. In this case the equalities of Theorem 27.14 are
true for $k=\mu \mathbf{a}$ with $1 \leq \mu \in 2^{-1} \mathbf{Z}$; if $\mu=1 / 2$, however, we need the residues of $E_{\mu \mathbf{a}}^{n, 0}(z, s)$ at $s=3 / 4$. For details, the reader is referred to [S85b]. Theorems 27.15 and 27.16 include the case in which $n=1$ and $r=0$; we have to assume $\mu \geq 3 / 2$ if $F \neq \mathbf{Q}$ and $\mu \geq 5 / 2$ if $F=\mathbf{Q}$. In fact, we can show that the orthogonal complement of $\mathcal{S}_{\mu \mathbf{a}}^{1}$ in $\mathcal{M}_{\mu \mathbf{a}}^{1}$ is spanned by $\mathbf{Q}$-rational elements for every $\mu \geq 1 / 2$. The proof is given in [ $\mathbf{S 8 7 b}$, Theorem 9.1]; the case in which $\mu=3 / 2$ and $F=\mathbf{Q}$ is complicated; we have to invoke the results of Pei $[\mathrm{P}]$ and Miyake $[\mathrm{M}]$.
(B) The case $\mu=\lambda(N)$ is not included in Theorem 27.16. In fact, in Corollaries 28.12 and 29.8 below we shall prove that $E_{k}^{n, r}(z, f, \Gamma)$ is $\overline{\mathbf{Q}}$-rational for every $\overline{\mathbf{Q}}$ rational $f$, if $m=\mu \mathbf{a}, \mu=\lambda(N), \mu>(3 r / 2)+1$ in Case SP, $\mu>3 r$ in Case UT, and $\mu>2 n$ in Case UB.
(C) The series defining $E_{k}^{\psi, \varphi_{r}}(z, s ; f, \Gamma)$ is convergent for $\operatorname{Re}(s)>(n+r+1) / 2$ in Case SP, for $\operatorname{Re}(s)>n+r$ in Case UT, and for $\operatorname{Re}(s)>l+r+\operatorname{dim}(Z)$ in Case UB. In these cases the value $E_{k}^{\psi, \varphi_{r}}(z, \mu / 2 ; f, \Gamma)$ is clearly holomorphic, and we do not need Lemma 27.10 for the proof of our lemmas and theorems after that lemma. Though the results even in such convergent cases are by no means trivial, it may be emphasized that one of the main points of our treatment in this section is in the fact that the above theorems are valid beyond the range of convergence.

## 28. Main theorems on arithmeticity in Cases SP and UT

28.1. Throughout this section we put $d=[F: \mathbf{Q}]$ and denote by $\Phi$ the Galois closure of $K$ in $\mathbf{C}$ over $\mathbf{Q}$; also we fix a weight $k \in \mathbf{Z}^{\mathbf{b}}$ and put $m=k$ in Case SP and $m=\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$ in Cases UT and UB; $k$ may be half-integral in Case SP. Let $G=S p(n, F)$ in Case SP, $G=U\left(\eta_{n}\right)$ in Case UT, and $G=G^{\varphi}$ with $(V, \varphi)$ as in Section 26 in Case UB. Given a congruence subgroup $\Gamma$ of $G$ and $f, g \in C^{\infty}(\mathcal{H}, \mathbf{C})$ such that $f \|_{k} \gamma=f$ and $g \|_{k} \gamma=g$ for every $\gamma \in \Gamma$, we define the inner product $\langle g, f\rangle$ by

$$
\begin{equation*}
\langle g, f\rangle=\operatorname{vol}(\mathfrak{D})^{-1} \int_{\mathfrak{D}} \overline{g(z)} f(z) \delta(z)^{m} \mathbf{d} z, \quad \operatorname{vol}(\mathfrak{D})=\int_{\mathfrak{D}} \mathbf{d} z, \quad \mathfrak{D}=\Gamma \backslash \mathcal{H} \tag{28.1}
\end{equation*}
$$

whenever the integral is convergent, where $\mathbf{d} z=\prod_{v \in \mathbf{a}} \mathbf{d} z_{v}$ with $\mathbf{d} z_{v}$ of Lemma 3.4 in Cases SP and UT; $\mathbf{d} z$ on $\mathfrak{Z}^{\varphi}$ is given by [S97, (10.9.1)]. This is essentially the same as that of (12.35a). Here, as well as in (28.2) below, we assume that $G_{\mathrm{a}}$ is not compact. For compact $G_{\mathbf{a}}$, inner products are defined in [S97, (10.9.5)], but we shall not employ them in the present book.

Let $C$ be an open subgroup of $G_{\mathbf{A}}$ such that $C \cap G_{\mathrm{h}}$ is compact. Take a finite subset $\mathcal{B}$ of $G_{\mathbf{h}}$ so that $G_{\mathbf{A}}=\bigsqcup_{p \in \mathcal{B}} G p C$. Let $W$ be a subfield of $\mathbf{C}$ such that $\Phi \mathbf{Q}_{\mathrm{ab}} \subset W$ in Cases SP and UT and $\overline{\mathbf{Q}} \subset W$ in Case UB. Let $\left(g_{p}\right)_{p \in \mathcal{B}} \leftrightarrow \mathbf{g} \in$ $\mathcal{M}_{k}(C)$ in the sense of $\S 20.1$ or $\S 26.10$ (or [S97, $\left.\S 10.7\right]$ ). We call $\mathbf{g} W$-rational if $g_{p} \in \mathcal{S}_{k}(W)$ for every $p \in \mathcal{B}$. If $q \in G_{\mathbf{h}} \cap G p C$, then $q=\alpha p u$ with $\alpha \in G$ and $u \in C$, so that $g_{q}=g_{p} \|_{k} \alpha^{-1}$. Therefore $g_{q} \in \mathcal{M}_{k}(W)$ by Theorems 9.13 (3), 10.7 (6), and Proposition 11.13. Thus the $W$-rationality is independent of the choice of $\mathcal{B}$. (In Case SP we can always take $\mathcal{B}=\{1\}$, and hence we can speak of the rationality of $g_{1}$ over any field, but we shall not employ it in this section.) Let $\mathcal{S}_{k}(C, W)$ denote the set of all $W$-rational elements of $\mathcal{S}_{k}(C)$. By Theorem 10.8 and Proposition 26.8 (2) we have $\mathcal{S}_{k}(C)=\mathcal{S}_{k}(C, W) \otimes_{W}$ C. In Cases SP and UT we easily see that $\mathbf{g} \in \mathcal{S}_{k}(C, W)$ if and only if $c_{\mathbf{g}}(\tau, q)$ of Proposition 20.2 belongs to $W$ for every $(\tau, q)$.

Now given also $\left(f_{b}\right)_{b \in \mathcal{B}} \leftrightarrow \mathbf{f} \in \mathcal{S}_{k}(C)$, we put

$$
\begin{equation*}
\langle\mathbf{g}, \mathbf{f}\rangle=\#(\mathcal{B})^{-1} \sum_{b \in \mathcal{B}}\left\langle g_{b}, f_{b}\right\rangle \quad \text { if } G_{\mathbf{a}} \text { is not compact. } \tag{28.2}
\end{equation*}
$$

This is well-defined independently of the choice of $C$ and $\mathcal{B}$; see [S97, §10.9].
In the following lemma we assume that $\varphi$ is isotropic in Case UB, since if $\varphi$ is anisotropic, then Proposition 15.7 guarantees the map $\mathfrak{p}$ with no condition on $k$ and $p$.
28.2. Lemma. Let $k$ be a weight and let $0 \leq p \in \mathbf{Z}^{\mathbf{a}}$. Suppose that for every $v \in \mathbf{a}$ we have $k_{v}>n+p_{v}$ in Case SP, $m_{v} \geq 2 n+p_{v}$ in Case UT, and $m_{v} \geq$ $\operatorname{dim}(V)+p_{v}$ in Case UB; suppose also that (27.12) is satisfied if $m=\mu \mathbf{a}$ with $\mu \in 2^{-1} \mathbf{Z}$. Let $D=\Phi \mathbf{Q}_{\mathrm{ab}}$ if $m \notin 2^{-1} \mathbf{Z}^{\mathbf{a}}$ in Cases SP and UT; let $D=\overline{\mathbf{Q}}$ otherwise. Then there exists a $\mathbf{C}$-linear map $\mathfrak{p}: \mathcal{N}_{k}^{p} \rightarrow \mathcal{S}_{k}$ with the following properties:
(1) $\langle g, f\rangle=\langle\mathfrak{p}(g), f\rangle$ for every $f \in \mathcal{S}_{k}$ and every $g \in \mathcal{N}_{k}^{p}$.
(2) $\mathfrak{p}(g)^{\sigma}=\mathfrak{p}\left(g^{\sigma}\right)$ for every $\sigma \in \operatorname{Aut}(\mathbf{C} / D)$ and every $g \in \mathcal{N}_{k}^{p}$.

Proof. Take the operator $\mathfrak{A}$ of Proposition 15.3. Given $g \in \mathcal{N}_{k}^{p}$, put $h=\mathfrak{A} g$. We have $h \in \mathcal{M}_{k}$ by that proposition, and $\langle g, f\rangle=\langle\mathfrak{A} g, f\rangle$ for every $f \in \mathcal{S}_{k}$ by Corollary 15.4. Since $\mathfrak{A}$ is a Q-rational polynomial of the operators $L_{\omega, v}$ for $v \in \mathbf{a}$ with $\omega(x)=\operatorname{det}(x)^{k}$, we see that $\mathfrak{A}\left(g^{\sigma}\right)=(\mathfrak{A} g)^{\sigma}$ for every $\sigma \in \operatorname{Aut}(\mathbf{C} / D)$ by Theorems 14.9 and 14.12. If $m \notin 2^{-1} \mathbf{Z}^{\mathbf{a}}$ then $\mathcal{S}_{k}=\mathcal{M}_{k}$, and so we have properties (1) and (2) with $\mathfrak{A}$ as $\mathfrak{p}$. If $m=\mu \mathbf{a}$ with $\mu \in 2^{-1} \mathbf{Z}$, then take $\mathfrak{q}$ of Theorem 27.14 and put $\mathfrak{p}=\mathfrak{q} \mathfrak{A}$. Then we obtain (1) and (2) from that theorem.
28.3. We now consider only Cases SP and UT, and take our setting to be that of Section 22; Case UB will be treated in Section 29. We take $C$ in the form (19.1) with $\mathfrak{e}=\mathfrak{c}$ and a Hecke eigenform $\mathbf{f} \in \mathcal{S}_{k}(C)$ as in $\S 20.6$ with trivial $\psi$; we also take $\kappa=0$ in (22.3b). Then the Euler $\mathfrak{p}$-factor of $\mathcal{Z}(s, \mathbf{f}, \chi)$ is 1 for $\mathfrak{p} \mid \mathfrak{c}$. We note that the square of $\operatorname{vol}(X)$ of (22.6) is a rational number by [S97, (18.9.3) and (18.9.4)]. Therefore from (22.9) and (22.18b) we obtain

$$
\begin{equation*}
\Gamma((s)) D\left(u s+s_{0}, \mathbf{f}, \chi\right)=\operatorname{vol}(\mathfrak{D}) \operatorname{det}(\tau)^{h+s \mathbf{a}} \sum_{p \in \mathcal{A}} a_{p} b_{p}^{u s}\left\langle g_{p} E\left(\bar{s}+\lambda_{n}\right), f_{p}\right\rangle \tag{28.3}
\end{equation*}
$$

with $a_{p} \in \overline{\mathbf{Q}}, 0<b_{p} \in \mathbf{Q}, E(s)=E\left(z, s ; m-m^{\prime}, 0, \Gamma\right)$, and a certain finite subset $\mathcal{A}$ of $G_{\mathbf{h}} ; a_{p} \in \mathbf{Q}_{\mathrm{ab}}$ in Case SP; $u=2, \lambda_{n}=(n+1) / 2$, and $s_{0}=(3 n / 2)+1$ in Case SP; $u=1, \lambda_{n}=n$, and $s_{0}=3 n / 2$ in Case UT; $h$ is defined by (22.4a) with $\kappa=0$.

Now $\operatorname{vol}(\mathfrak{D})$ is the Euler-Poincaré characteristic of $\Gamma \backslash \mathcal{H}$ times a constant that depends only on the choice of the measure of $\mathcal{H}$. (This is a well-known principle valid for any arithmetic quotient of a hermitian symmetric space. If the quotient is compact, it follows from the classical generalization of the Gauss-Bonnet formula; the noncompact case was proved in [Ha].) In the present case, the constant is $\pi^{d_{0}}$ times a rational number, where $d_{0}$ is the complex dimension of $\mathcal{H}$. (The constant depends on the type of $G_{v}$. In the symplectic case, the rationality follows, for example, from [Si, II, p.279, Theorem 11], which gives $\operatorname{vol}(\mathcal{D})$ when $\Gamma=S p(n, \mathbf{Z})$. In the unitary case it follows from the formula for $\operatorname{vol}(\mathcal{D})$ in [S97, Proposition 24.9], in which we can take $\operatorname{det}\left(\theta_{v}\right)=1$ for every $v \in \mathbf{a}$.) Therefore, using the notation of (22.19) and putting $\sigma=u s+s_{0}$, we obtain

$$
\begin{equation*}
c_{\mathbf{f}}(\tau, r) \Gamma((s)) \mathcal{Z}(\sigma, \mathbf{f}, \chi) \tag{28.4}
\end{equation*}
$$

$$
=\operatorname{det}(\tau)^{h+s \mathbf{a}} P(\sigma) \Lambda(\sigma) \pi^{d_{0}} \sum_{p \in \mathcal{A}} a_{p}^{\prime} b_{p}^{u s}\left\langle g_{p} E\left(\bar{s}+\lambda_{n}\right), f_{p}\right\rangle
$$

with $a_{p}^{\prime} \in \overline{\mathbf{Q}} ; a_{p}^{\prime} \in \mathbf{Q}_{\mathrm{ab}}$ in Case SP. Here $\Lambda, P, a_{p}^{\prime}$, and $b_{p}$ depend only on ( $\tau, r$ ). Strictly speaking, (28.4) holds for a Hecke eigenform $\mathbf{f}$ with a special property relative to the pair ( $\tau, r$ ) as explained in Theorem 20.9; also, given $\mathbf{f}$, we can always find ( $\tau, r$ ) with which (28.4) holds. Now, for an arbitrary eigenform $\mathbf{f}$ and arbitrary $(\tau, r)$ we have

$$
\begin{align*}
& \Gamma((s)) \mathcal{Z}(\sigma, \mathbf{f}, \chi) \sum_{L<M \in \mathcal{L}_{\tau}} \mu(M / L) \chi^{*}\left(\operatorname{det}\left(q^{*} \widehat{y}\right) \mathfrak{r}\right) N\left(\operatorname{det}\left(q^{*} \widehat{y}\right) \mathfrak{r}\right)^{-\sigma} c_{\mathbf{f}}(\tau, y)  \tag{28.5}\\
& =\operatorname{det}(\tau)^{h+s \mathbf{a}} P(\sigma) \Lambda(\sigma) \pi^{d_{0}} \sum_{p \in \mathcal{A}} a_{p}^{\prime} b_{p}^{u s}\left\langle g_{p} E\left(\bar{s}+\lambda_{n}\right), f_{p}\right\rangle
\end{align*}
$$

with the same $\Lambda, P, a_{p}^{\prime}$, and $b_{p}$, where the sum on the left-hand side is as in Theorem 20.7. Indeed, we obtained (22.19) from (20.19), and (20.19) is a special case of the equality of Theorem 20.7. Therefore (28.5) can be obtained by employing that equality instead of (20.19). The function $\Lambda$ is obtained from $\mathfrak{L}_{0} \prod_{v \nmid c} h_{v}$ of Theorem 20.9 as explained in $\S 22.9 ; \mathfrak{L}_{0}$ and $h_{v}$ for $k \notin \mathbf{Z}^{\mathbf{a}}$ are given in Theorem 21.4. Thus we obtain

$$
\begin{aligned}
& \Lambda(\sigma)=L_{\mathfrak{c}}\left(\sigma-\frac{n}{2}, \chi \rho_{\tau}\right) \prod_{i=1}^{n / 2} L_{\mathfrak{c}}\left(2 \sigma-2 n-2+2 i, \chi^{2}\right) \quad\left(\text { Case SP, } k \in \mathbf{Z}^{\mathbf{a}}, n \in 2 \mathbf{Z}\right), \\
& \Lambda(\sigma)=\prod_{i=1}^{(n+1) / 2} L_{\mathfrak{c}}\left(2 \sigma-2 n-2+2 i, \chi^{2}\right) \quad\left(\text { Case SP, } k \in \mathbf{Z}^{\mathbf{a}}, n \notin 2 \mathbf{Z}\right), \\
& \Lambda(\sigma)=\prod_{i=1}^{n / 2} L_{\mathfrak{c}}\left(2 \sigma-2 n-1+2 i, \chi^{2}\right) \quad\left(\text { Case SP, } k \notin \mathbf{Z}^{\mathbf{a}}, n \in 2 \mathbf{Z}\right), \\
& \Lambda(\sigma)=L_{\mathfrak{c}}\left(\sigma-\frac{n}{2}, \chi \rho_{\tau}\right) \prod_{i=1}^{(n-1) / 2} L_{\mathfrak{c}}\left(2 \sigma-2 n-1+2 i, \chi^{2}\right) \quad\left(\text { Case SP, } k \notin \mathbf{Z}^{\mathbf{a}}, n \notin 2 \mathbf{Z}\right), \\
& \Lambda(\sigma)=\prod_{i=1}^{n} L_{\mathfrak{c}}\left(2 \sigma-n+1-i, \chi_{1} \theta^{n+i-1}\right) \quad(\text { Case UT). }
\end{aligned}
$$

Here $\rho_{\tau}$ is the Hecke character of $F$ given in Lemma 20.5 and Theorem 21.4; $\chi_{1}$ is the restriction of $\chi$ to $F$.
28.4. Lemma. The notation being as above, let $h_{p}$ be an element of $\mathcal{S}_{k}(D)$ given for each $p \in \mathcal{A}$, where $D$ is a subfield of $\mathbf{C}$ containing $\Phi \mathbf{Q}_{\mathbf{a b}}$. Then there exists an element $\mathbf{q}$ of $\mathcal{S}_{k}(C, D)$ independent of $\mathbf{f}$ such that $\sum_{p \in \mathcal{A}}\left\langle h_{p}, f_{p}\right\rangle=\langle\mathbf{q}, \mathbf{f}\rangle$.

Proof. Recall that $f_{p} \in \mathcal{S}_{k}\left(\Gamma^{p}\right)$ with $\Gamma^{p}=G \cap p C p^{-1}$. Take $\Gamma \subset \bigcap_{p \in \mathcal{A}} \Gamma^{p}$ so that $h_{p} \in \mathcal{S}_{k}(\Gamma)$ for every $p \in \mathcal{A}$; let $\Gamma^{p}=\bigsqcup_{\xi \in X_{p}} \Gamma \xi$. Put then $h_{p}^{\prime}=\#\left(X_{p}\right)^{-1} \sum_{\xi \in X_{p}} h_{p} \| \xi$. Then $h_{p}^{\prime} \in \mathcal{S}_{k}\left(\Gamma^{p}, D\right)$ and $\left\langle h_{p}, f_{p}\right\rangle=\left\langle h_{p}^{\prime}, f_{p}\right\rangle$. For each $b \in \mathcal{B}$ let $\mathcal{A}_{b}=\{p \in \mathcal{A} \mid p \in$ $G b C\}$ and $q_{b}=\#(\mathcal{B}) \sum_{p \in \mathcal{A}_{b}} h_{p}^{\prime} \| \alpha_{p}$, where $\alpha_{p} \in G$ is chosen so that $p \in \alpha_{p} b C$. Define $\mathbf{q} \in \mathcal{S}_{k}(C, D)$ by $\mathbf{q} \leftrightarrow\left(q_{b}\right)_{b \in \mathcal{B}}$. Then $\left\langle q_{b}, f_{b}\right\rangle=\#(\mathcal{B}) \sum_{p \in \mathcal{A}_{b}}\left\langle h_{p}^{\prime} \| \alpha_{p}, f_{b}\right\rangle$ $=\#(\mathcal{B}) \sum_{p \in \mathcal{A}_{b}}\left\langle h_{p}^{\prime}, f_{p}\right\rangle$, and hence we obtain the desired conclusion.
28.5. Theorem (Cases SP and UT). Let $\{\lambda(\mathfrak{a})\}$ be a system of eigenvalues on $\mathcal{S}_{k}(C)$ in the sense that $\mathbf{f}_{0} \mid T(\mathfrak{a})=\lambda(\mathfrak{a}) \mathbf{f}_{0}$ for every $\mathfrak{a}$ with some $\mathbf{f}_{0} \in \mathcal{S}_{k}(C), \neq 0$ (see §20.6). Let $\Psi$ be the field generated by the $\lambda(\mathfrak{a})$ over $\Phi \mathbf{Q}_{\mathrm{ab}} ;$ put

$$
\begin{equation*}
\mathcal{V}=\left\{\mathbf{f} \in \mathcal{S}_{k}(C)|\mathbf{f}| T(\mathfrak{a})=\lambda(\mathfrak{a}) \mathbf{f} \quad \text { for every } \mathfrak{a}\right\} \tag{28.6}
\end{equation*}
$$

$$
\begin{equation*}
\mathcal{V}(\Psi)=\mathcal{V} \cap \mathcal{S}_{k}(C, \Psi), \quad \mathcal{V}(\overline{\mathbf{Q}})=\mathcal{V} \cap \mathcal{S}_{k}(C, \overline{\mathbf{Q}}) \tag{28.6a}
\end{equation*}
$$

Then $\Psi$ is stable under complex conjugation and $\mathcal{V}=\mathcal{V}(\Psi) \otimes_{\Psi}$ C. Moreover, put $m_{0}=\operatorname{Min}_{v \in \mathbf{a}} m_{v}$ and assume that the following condition is satisfied:
(28.7) $m_{0}>(3 n / 2)+1$ in Case SP and $m_{0}>3 n$ in Case UT.

Then

$$
\left\langle\mathbf{g}, \mathbf{g}^{\prime}\right\rangle /\langle\mathbf{f}, \mathbf{f}\rangle \in \overline{\mathbf{Q}} \text { if } \mathbf{g}, \mathbf{g}^{\prime} \in \mathcal{V}(\overline{\mathbf{Q}}) \text { and } 0 \neq \mathbf{f} \in \mathcal{V}(\overline{\mathbf{Q}})
$$

Proof. That $\Psi$ is stable under complex conjugation follows from Lemma 23.15. By Lemma $20.12(3), \mathcal{S}_{k}(C)$ is spanned by eigenforms of the operators $T(\mathfrak{a})$, which are normal. From [S97, (11.9.1)], Theorem 9.13 (3), (21.4), and Theorem 10.7 (6) we see that the $T(\mathfrak{a}) \operatorname{map} \mathcal{S}_{k}(C, \Psi)$ into itself, and so they generate a ring of semisimple $\Psi$-linear transformations on $\mathcal{S}_{k}(C, \Psi)$. Therefore we have $\mathcal{V}=\mathcal{V}(\Psi) \otimes_{\Psi} \mathbf{C}$ and $\mathcal{S}_{k}(C, \Psi)=\mathcal{V}(\Psi) \oplus \mathcal{U}$ with a vector subspace $\mathcal{U}$ over $\Psi$ stable under the $T(\mathfrak{a})$. Each eigenform in $\mathcal{U} \otimes_{\Psi} \mathbf{C}$, being not contained in $\mathcal{V}$, must be orthogonal to $\mathcal{V}$. Thus $\mathcal{U}$ is orthogonal to $\mathcal{V}$.

To prove the main part of our theorem, put $\mathcal{Z}(s)=\mathcal{Z}(s, \mathbf{f}, \chi)$. This depends only on the $\lambda(\mathfrak{a})$. We first consider Case SP. Define $\mu \in \mathbf{Z}^{\mathfrak{a}}$ so that $0 \leq \mu_{v} \leq 1$ and $m_{0}-k_{v}+\mu_{v} \in 2 \mathbf{Z}$. Put $l=\mu+(n / 2) \mathbf{a}, t^{\prime}=\mu-[k], \nu=m_{0}-(n / 2)$, and $\sigma_{0}=m_{0}$. By [S97, Lemma 11.14 (3)] we can find a Hecke character $\chi$ of $F$ such that $\chi_{\mathbf{a}}(x)=x_{\mathbf{a}}^{-t^{\prime}}\left|x_{\mathbf{a}}\right| t^{t^{\prime}}$. We consider (28.4) by employing $\mathbf{g}$ defined as in §A5.5 with such $\chi$ and $\mu$; then $g_{p} \in \mathcal{M}_{k}\left(\Phi \mathbf{Q}_{\mathrm{ab}}\right)$. Observe that $n+1<\nu \leq k_{v}-l_{v}$ and $k_{v}-l_{v}-\nu \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$. Evaluate (28.4) at $s=(\nu-n-1) / 2$ (which means that $\left.\sigma=\sigma_{0}\right)$. Putting $Q=\Gamma(((\nu-n-1) / 2))^{-1} \Lambda\left(\sigma_{0}\right)$, we thus obtain

$$
\begin{equation*}
c_{\mathbf{f}}(\tau, r) \mathcal{Z}\left(\sigma_{0}\right)=\pi^{d_{0}} \operatorname{det}(\tau)^{h+s \mathbf{a}} P_{\tau, r}\left(\sigma_{0}\right) Q \sum_{p \in \mathcal{A}}\left\langle e_{p} g_{p} E(\nu / 2), f_{p}\right\rangle \tag{28.8}
\end{equation*}
$$

with $e_{p} \in \mathbf{Q}_{\mathrm{ab}}$; we write $P_{\tau, r}$ for $P$ in order to emphasize its dependence on $(\tau, r)$. From (22.4a) we see that $s+h_{v}-m_{0}+(n+1) / 2 \in \mathbf{Z}$ for every $v \in \mathbf{a}$, so that $\operatorname{det}(\tau)^{h+s \mathbf{a}} \in \Phi \mathbf{Q}_{\mathrm{ab}}$; also we can easily verify that $Q \neq 0$. By (17.21) and Theorem 17.7 (i) we can put $E(\nu / 2)=\Delta_{\nu \mathbf{a}}^{q} y$ with $y \in \mathcal{M}_{\nu \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ and $q=(k-l-\nu \mathbf{a}) / 2$. Therefore, by Lemma 15.8, there exists an element $h_{p}^{\prime}$ of $\mathcal{M}_{k}\left(\Phi \mathbf{Q}_{\mathrm{ab}}\right)$ such that $\left\langle e_{p} g_{p} E(\nu / 2), f_{p}\right\rangle=\pi^{n|q|}\left\langle h_{p}^{\prime}, f_{p}\right\rangle$. Put $h_{p}=h_{p}^{\prime}$ if $k \notin 2^{-1} \mathbf{Z}^{\mathbf{a}}$; otherwise put $h_{p}=\mathfrak{q}\left(h_{p}^{\prime}\right)$ with the map $\mathfrak{q}$ of Theorem 27.14. (Recall that $\mathcal{M}_{k}=\mathcal{S}_{k}$ if $k \notin 2^{-1} \mathbf{Z}^{\mathbf{a}}$ by [S97, Proposition 10.6 (3)].) Then $h_{p} \in \mathcal{S}_{k}(\overline{\mathbf{Q}})$ and $\left\langle h_{p}^{\prime}, f_{p}\right\rangle=\left\langle h_{p}, f_{p}\right\rangle$ since $\left\langle\mathcal{S}_{k}^{\psi}, \mathcal{E}_{k}^{\psi}\right\rangle=0$ if $k \in 2^{-1} \mathbf{Z}^{\mathbf{a}}$. By Lemma 28.4 we find an element $\mathbf{q}$ of $\mathcal{S}_{k}(C, \overline{\mathbf{Q}})$ independent of $\mathbf{f}$ such that $\sum_{p \in \mathcal{A}}\left\langle h_{p}, f_{p}\right\rangle=\langle\mathbf{q}, \mathbf{f}\rangle$. Clearly $\mathcal{S}_{k}(C, \overline{\mathbf{Q}})=\mathcal{V}(\overline{\mathbf{Q}}) \oplus$ $\left(\mathcal{U} \otimes_{\Psi} \overline{\mathbf{Q}}\right)$, and $\mathcal{V}(\overline{\mathbf{Q}})$ is orthogonal to $\mathcal{U} \otimes_{\Psi} \overline{\mathbf{Q}}$. Let $\mathbf{q}_{\tau, r}$ be the projection of $\mathbf{q}$ to $\mathcal{V}(\overline{\mathbf{Q}})$ with respect to this direct sum decomposition of $\mathcal{S}_{k}(C, \overline{\mathbf{Q}})$. Once $\mu, \chi$, and $(\tau, r)$ are fixed, then $g_{p}, P_{\tau, r}$, and $E$ are independent of $\mathbf{f}$. Thus given $0 \neq \mathbf{f} \in \mathcal{V}$, there exists $(\tau, r)$ such that $c_{\mathbf{f}}(\tau, r) \neq 0$ and

$$
\begin{equation*}
c_{\mathbf{f}}(\tau, r) \mathcal{Z}\left(\sigma_{0}\right)=\pi^{d_{0}+n|q|} \operatorname{det}(\tau)^{h+s \mathbf{a}} P_{\tau, r}\left(\sigma_{0}\right) Q\left\langle\mathbf{q}_{\tau, r}, \mathbf{f}\right\rangle \tag{28.9}
\end{equation*}
$$

Let $\mathfrak{S}$ denote the set of all pairs $(\tau, r)$ for which such an $\mathbf{f}$ exists. Since $\sigma_{0}>3 n / 2+$ $1, \mathcal{Z}\left(\sigma_{0}\right) \neq 0$ by Theorem 20.13. Therefore, given $\mathbf{f} \in \mathcal{V}$, we have $\left\langle\mathbf{q}_{\tau, r}, \mathbf{f}\right\rangle \neq 0$ for some $(\tau, r) \in \mathfrak{S}$, which implies that the $\mathbf{q}_{\tau, r}$ for all $(\tau, r) \in \mathfrak{S}$ generate $\mathcal{V}$ over $\mathbf{C}$, and hence they generate $\mathcal{V}(\overline{\mathbf{Q}})$ over $\overline{\mathbf{Q}}$. Also we see that $P_{\tau . r}\left(\sigma_{0}\right) \neq 0$ for
every $(\tau, r) \in \mathfrak{S}$. Now, given $\mathbf{f} \in \mathcal{V}(\overline{\mathbf{Q}})$ and an arbitrary $(\tau, r) \in \mathfrak{S}$, we may not be able to use (28.8), but evaluating (28.5) at $s=(\nu-n-1) / 2$, we find that $\left\langle\mathbf{q}_{\tau, r}, \mathbf{f}\right\rangle \in \pi^{-d_{0}-n|q|} Q^{-1} \mathcal{Z}\left(\sigma_{0}\right) \overline{\mathbf{Q}}$. Therefore $\langle\mathbf{g}, \mathbf{f}\rangle \in \pi^{-d_{0}-n|q|} Q^{-1} \mathcal{Z}\left(\sigma_{0}\right) \overline{\mathbf{Q}}$ for every $\mathbf{g}, \mathbf{f} \in \mathcal{V}(\overline{\mathbf{Q}})$. The main assertion of our theorem immediately follows from this fact.

Next let us treat Case UT. Define $\mu \in \mathbf{Z}^{\mathbf{b}}$ as follows: $\mu_{v \rho}=0$ and $\mu_{v}=m_{v}-m_{0}$ for every $v \in \mathbf{a}$; put $t^{\prime}=\left(\mu_{v}-k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$. Then (22.15a, b) are satisfied. By [S97, Lemma 11.14 (3)] we can find a Hecke character $\chi$ of $K$ such that $\chi_{\mathbf{a}}(x)=$ $x_{\mathbf{a}}^{-t^{\prime}}\left|x_{\mathbf{a}}\right|^{t^{\prime}}$. Put $\nu=m_{0}-n$ and $\sigma_{0}=m_{0} / 2$. Define $\mathbf{g}$ as in $\S A 5.5$ and $m^{\prime}$ as in Proposition 22,2; then $m_{v}^{\prime}=\mu_{v}+n$ for $v \in \mathbf{a}$, and hence $2 n<\nu=m_{v}-m_{v}^{\prime}$ for every $v \in \mathbf{a} ; g_{p} \in \mathcal{M}_{\mu+n \mathbf{a}}(\overline{\mathbf{Q}})$, since $\chi^{*}$ is algebraic-valued by Lemma 17.11. Evaluate (28.4) at $s=\left(m_{0}-3 n\right) / 2$. Then $\sigma$ and $E\left(\bar{s}+\lambda_{n}\right)$ become $\sigma_{0}$ and $E(\nu / 2)$, and $m-m^{\prime}=\nu \mathbf{a} ; \operatorname{det}(\tau)^{h+s \mathbf{a}} \in \Phi$. By Theorem 17.7 (i) we find that $E(\nu / 2) \in \mathcal{M}_{\nu \mathbf{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ so that $e_{p} g_{p} E(\nu / 2) \in \mathcal{M}_{k}(\overline{\mathbf{Q}})$. Then we can repeat our argument in Case SP, by putting $h_{p}=\mathfrak{q}\left(e_{p} g_{p} E(\nu / 2)\right)$ without employing $\Delta_{\nu \mathrm{a}}^{p}$.
28.6. Corollary. The notation being the same as in (28.2) and Theorem 28.5, let $0 \neq \mathbf{f} \leftrightarrow\left(f_{b}\right)_{b \in \mathcal{B}} \in \mathcal{V}(\overline{\mathbf{Q}})$. Then $\left\langle g, f_{a}\right\rangle \in\langle\mathbf{f}, \mathbf{f}\rangle \overline{\mathbf{Q}}$ for every $g \in \mathcal{M}_{k}(\overline{\mathbf{Q}})$ and every $a \in \mathcal{B}$.

Proof. We may assume that $g \in \mathcal{S}_{k}(\overline{\mathbf{Q}})$. Indeed, $\mathcal{M}_{k} \neq \mathcal{S}_{k}$ only if $m=\mu \mathbf{a}$ with $\mu \in 2^{-1} \mathbf{Z}$, in which case we put $g^{\prime}=\mathfrak{q}(g)$ with the map $\mathfrak{q}$ of Theorem 27.14. Then $g^{\prime} \in \mathcal{S}_{k}(\overline{\mathbf{Q}})$ and $\langle g, f\rangle=\left\langle g^{\prime}, f\right\rangle$ for every $f \in \mathcal{S}_{k}$, so that it is sufficient to treat the case $g \in \mathcal{S}_{k}(\overline{\mathbf{Q}})$. Given $g \in \mathcal{S}_{k}(\overline{\mathbf{Q}})$, we may assume, changing $C$ for its suitable subgroup, that $g \in \mathcal{S}_{k}\left(\Gamma^{a}\right)$, where $\Gamma^{a}=G \cap a C a^{-1}$. Fixing $a \in \mathcal{B}$, define $\mathbf{g} \leftrightarrow\left(g_{b}\right)_{b \in \mathcal{B}} \in \mathcal{S}_{k}(C, \overline{\mathbf{Q}})$ so that $g_{a}=g$ and $g_{b}=0$ for $a \neq b \in \mathcal{B}$. We have $\mathcal{S}_{k}(C, \overline{\mathbf{Q}})=\mathcal{V}(\overline{\mathbf{Q}}) \oplus\left(\mathcal{U} \otimes_{\mathcal{L}} \overline{\mathbf{Q}}\right)$ with $\mathcal{U}$ as in the proof of Theorem 28.5. Let $\mathbf{g}^{\prime}$ be the projection of $\mathbf{g}$ to $\mathcal{V}(\overline{\mathbf{Q}})$ with respect to that decomposition. Then $\#(\mathcal{B})^{-1}\left\langle g, f_{a}\right\rangle=\langle\mathbf{g}, \mathbf{f}\rangle=\left\langle\mathbf{g}^{\prime}, \mathbf{f}\right\rangle$, which belongs to $\langle\mathbf{f}, \mathbf{f}\rangle \overline{\mathbf{Q}}$ by Theorem 28.5.
28.7. Lemma. In the setting of Sections 24 and 25, let $R \in \mathcal{M}_{\nu \mathbf{a}}^{n+r}(\Xi)$ with a subfield $\Xi$ of $\mathbf{C}$ containing $\mathbf{Q}_{\mathrm{ab}}$ and the Galois closure of $K$ over $\mathbf{Q}$, and let $S_{0}=D_{e, e^{\prime}} \Delta_{\nu \mathbf{a}}^{p} R$ and $S(z, w)=S_{0}(\operatorname{diag}[z, w])$ for $(z, w) \in \mathcal{H}^{n} \times \mathcal{H}^{r}$, where $D_{e, e^{\prime}}$, which we ignore if $n \neq r$, is as in §25.2, and $\Delta_{\nu \mathbf{a}}^{p}$ with $0 \leq p \in \mathbf{Z}^{\mathbf{a}}$ is as in Lemma 15.8 and the proof of Theorem 17.9. Let $L_{v}$ be the operator of (15.3) defined on $\mathcal{H}^{r}$ with $\omega(a, b)=\operatorname{det}(b)^{m}, m=\nu \mathbf{a}+2 p+e-e^{\prime}$. Then there exists an element $T(z, w)$ of $\sum_{v \in \mathbf{a}} \sum_{i=0}^{\infty} \Xi L_{v}^{i} S$ which is holomorphic in $w$ and such that $\langle S(z, w), f(w)\rangle=\langle T(z, w), f(w)\rangle$ for every $f \in \mathcal{S}_{m}^{r}$, provided $\nu \geq(n+r) / 2$ in Case SP and $\nu \geq n+r$ in Case UT.

This will be proven in §A8.12.
28.8. Theorem (Cases SP and UT). The notation being as in Theorem 28.5, let $0 \neq \mathbf{f} \in \mathcal{V}(\overline{\mathbf{Q}})$; let $\chi$ be a Hecke character of $K$ such that $\chi_{\mathbf{a}}(x)=x_{\mathbf{a}}^{\ell}\left|x_{\mathbf{a}}\right|^{-\ell}$ with $\ell \in \mathbf{Z}^{\mathbf{a}}$, and let $\sigma_{0} \in 2^{-1} \mathbf{Z}$. In addition to (28.7), assume the following condition:

Case SP: $2 n+1-k_{v}+\mu_{v} \leq \sigma_{0} \leq k_{v}-\mu_{v}$ where $\mu_{v}=0$ if $\left[k_{v}\right]-\ell_{v} \in 2 \mathbf{Z}$ and $\mu_{v}=1$ if $\left[k_{v}\right]-\ell_{v} \notin 2 \mathbf{Z} ; \sigma_{0}-k_{v}+\mu_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$ if $\sigma_{0}>n$ and $\sigma_{0}-1+k_{v}-\mu_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$ if $\sigma_{0} \leq n$.
Case UT: $4 n-\left(2 k_{v \rho}+\ell_{v}\right) \leq 2 \sigma_{0} \leq m_{v}-\left|k_{v}-k_{v \rho}-\ell_{v}\right|$ and $2 \sigma_{0}-\ell_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$.

Further exclude the following cases:
(A) Case SP: $\sigma_{0}=n+1, F=\mathbf{Q}$, and $\chi^{2}=1$;
(B) Case SP: $\sigma_{0}=n+(3 / 2), F=\mathbf{Q}, \chi^{2}=1$ and $[k]-\ell \in 2 \mathbf{Z}$;
(C) Case SP: $\sigma_{0}=0, \mathfrak{c}=\mathfrak{g}$, and $\chi=1$;
(D) Case SP: $0<\sigma_{0} \leq n, \mathfrak{c}=\mathfrak{g}, \chi^{2}=1$, and the conductor of $\chi$ is $\mathfrak{g}$;
(E) Case UT: $2 \sigma_{0}=2 n+1, F=\mathbf{Q}, \chi_{1}=\theta$, and $k_{v}-k_{v \rho}=\ell_{v}$;
(F) Case UT: $0 \leq 2 \sigma_{0}<2 n, \mathfrak{c}=\mathfrak{g}, \chi_{1}=\theta^{2 \sigma_{0}}$, and the conductor of $\chi$ is $\mathfrak{r}$.

Here $\chi_{1}$ is the restriction of $\chi$ to $F_{\mathbf{A}}^{\times}$and $\theta$ is the Hecke character of $F$ corresponding to $K / F$. Then

$$
\begin{equation*}
\mathcal{Z}\left(\sigma_{0}, \mathbf{f}, \chi\right) /\langle\mathbf{f}, \mathbf{f}\rangle \in \pi^{n|m|+d \varepsilon} \overline{\mathbf{Q}} \tag{28.10}
\end{equation*}
$$

where $d=[F: \mathbf{Q}],|m|=\sum_{v \in \mathbf{a}} m_{v}$, and

$$
\varepsilon= \begin{cases}(n+1) \sigma_{0}-n^{2}-n & \left(\text { Case SP, } k \in \mathbf{Z}^{\mathbf{a}} \text { and } \sigma_{0}>n\right) \\ n \sigma_{0}-n^{2} & \left(\text { Case SP, } k \notin \mathbf{Z}^{\mathbf{a}} \text { or } \sigma_{0} \leq n\right) \\ 2 n \sigma_{0}-2 n^{2}+n & \text { (Case UT). }\end{cases}
$$

Notice that $n|m|+d \varepsilon \in \mathbf{Z}$ in all cases. If $k \notin \mathbf{Z}^{\mathbf{a}}$, for example, the above condition on $\sigma_{0}$ shows that $\sigma_{0}+k_{v} \in \mathbf{Z}$, which implies that $n|m|+d \varepsilon \in \mathbf{Z}$.

Proof. There are two ways to prove this: the first one applies to the whole critical strip, and the second one only to the right half of the strip. However, the latter can cover certain cases to which the former does not apply. Let us begin with the first method, using the notation of $\S \S 25.4$ and 25.5 , in which the weight of $\mathbf{f}$ was written $h$ instead of $k$; at the end of the proof we shall reinstate $k$ as the weight of $\mathbf{f}$, and obtain the condition on $\sigma_{0}$ in terms of $k$ as stated above.

Cast UT (1st method). Define $d_{v}, e, e^{\prime}$, and $k$ as in $\S 25.4 ;$ put $m=\left(k_{v}+\right.$ $\left.k_{v \rho}\right)_{v \in \mathbf{a}}$ and $m^{\prime}=\left(h_{v}+h_{v \rho}\right)_{v \in \mathbf{a}}$. (At the end we must change $m^{\prime}$ into $m$.) We evaluate (25.8a) at $s=\sigma_{0}$; we have to change $\mathfrak{c}$, as we did in $\S 25.4$, so that the conductor of $\chi$ divides $\mathfrak{c}$. Put $\nu=2 \sigma_{0}$ and $H_{b, a}^{\prime}(\mathfrak{z}, s)=\Lambda_{\mathfrak{c}}^{2 n}(s, \chi) H_{b, a}(\mathfrak{z}, s)$ for $\mathfrak{z} \in \mathcal{H}^{2 n}$. We first treat the case $2 \sigma_{0}<2 n$. For the reason explained at the beginning of $\S 25.5, H_{b, a}^{\prime}$ is a function of type (17.24). Therefore, by (17.30), $H_{b, a}^{\prime}(\mathfrak{z}, \nu / 2)=$ $\Delta_{\mu \mathbf{a}}^{p} R_{1}$ with $R_{1}(\mathfrak{z})=D_{r}\left(\mathfrak{z}, \nu / 2, k^{\prime}, \chi, \mathfrak{c}\right)$, where $\mu=4 n-\nu, p=(m-\mu \mathbf{a}) / 2$, and $k^{\prime}$ is such that $\left(k_{v}^{\prime}+k_{v \rho}^{\prime}\right)_{v \in \mathbf{a}}=\mu \mathbf{a}$. By Theorem 17.12 (iii) we can put $R_{1}=\pi^{\alpha} R$ with $\alpha=d n(2 n+1)$ and $R \in \mathcal{M}_{\mu \mathbf{a}}(\overline{\mathbf{Q}})$. (Our $H^{\prime}$ is a function on $\mathcal{H}^{2 n}$, so that we have to take $2 n$ in place of $n$ in Theorem 17.7. For example, $\pi^{-2 n|p|} \Delta_{\mu \mathbf{a}}^{p} R$ is $\overline{\mathbf{Q}}$-rational.) We have to assume that $0 \leq m_{v}-\mu \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$. Since $m_{v}=h_{v}+h_{v \rho}-d_{v}=2 h_{v \rho}+\ell_{v}$, this means that

$$
\begin{equation*}
4 n-\left(2 h_{v \rho}+\ell_{v}\right) \leq \nu<2 n \quad \text { and } \quad \nu-\ell_{v} \in 2 \mathbf{Z} \quad \text { for every } v \in \mathbf{a} \tag{i}
\end{equation*}
$$

Case (F) of our theorem must be excluded. In (25.8a) we have $c_{m^{\prime}}\left(\mathbf{s}^{\prime}\right) \in \pi^{d_{0}} \mathbf{Q}^{\times}$ for the same reason as in (25.11); employing the formulas of Theorem 12.13 and Lemma 25.3, we see that $\Psi(\mathbf{s}) \neq 0$ if in addition to (i) we assume

$$
\begin{equation*}
\nu \leq m_{v}+2 d_{v} \quad \text { whenever } \quad d_{v}<0 \tag{ii}
\end{equation*}
$$

Define $S_{0}$ as in Lemma 28.7 for the present $R$. Then $\pi^{\alpha} S=\left[D_{e . e^{\prime}} H_{b . a}^{\prime}(z, \nu / 2)\right]^{\circ}$ and $\pi^{n\left|e^{\prime}-e-2 p\right|} S$ is a $\overline{\mathbf{Q}}$-rational function, to which Lemma 25.8 is applicable. Take $T(z, w)$ as in Lemma 28.7 and put $T_{a, b}(z, w)=T(z, \eta w) j_{\eta}^{h}(w)^{-1}$. For the same
reason as in Lemma 25.8, we can put $T_{a, b}(z, w)=\pi^{n\left|2 p+e-e^{\prime}\right|} \sum_{i} t_{a b i}(z) g_{a b i}(w)$ with $t_{a b i} \in \mathcal{N}_{h}^{p^{\prime}}(\overline{\mathbf{Q}})$ and $g_{a b i} \in \mathcal{M}_{h}(\overline{\mathbf{Q}})$. The integral over $\mathcal{D}_{a}$ of (25.8a) is a constant times $\left\langle f_{a}^{\prime}, S(z, \eta w) j_{\eta}^{h}(w)^{-1}\right\rangle$, where $f_{a}^{\prime}(w)=\overline{f_{a}\left(-w^{*}\right)}$. By the property of $T$ in Lemma 28.7, the last inner product equals $\left\langle f_{a}^{\prime}, T_{a, b}(z, w)\right\rangle$, which means that we can replace $\Lambda_{c}^{2 n} J_{b, a}^{\prime}$ at $s=\nu / 2$ in (25.8a) by $\pi^{\alpha} T_{a, b}\left(z,-w^{*}\right)$. Putting $g_{a b i}^{\prime}(w)=$ $\overline{g_{a b i}\left(-w^{*}\right)}$ and changing $t_{a b i}$ for its suitable scalar multiple, we find that

$$
\begin{equation*}
\mathcal{Z}\left(\sigma_{0}, \mathbf{f}, \chi\right) f_{b}(z)=\pi^{\gamma} \sum_{a, i}\left\langle g_{a b i}^{\prime}, f_{a}\right\rangle t_{a b i}^{\prime}(z) \tag{iii}
\end{equation*}
$$

with $\gamma=\alpha+n\left|2 p+e-e^{\prime}\right|$ and some $t_{a b i}^{\prime} \in \mathcal{N}_{h}^{p^{\prime}}(\overline{\mathbf{Q}})$. Take the Fourier expansion with respect to $z$, and compare nonzero Fourier coefficients. Then we find that $\mathcal{Z}\left(\sigma_{0}, \mathbf{f}, \chi\right)=\pi^{\gamma} \sum_{a}\left\langle g_{a}, f_{a}\right\rangle$ with some $g_{a} \in \mathcal{M}_{h}(\overline{\mathbf{Q}})$. Applying Corollary 28.6 to the right-hand side, we obtain (28.10). (If $\nu$ satisfies (i) and $d_{v} \geq 0$, then $2 n<2 h_{v \rho}+\ell_{v}=h_{v}+h_{v \rho}-\left|d_{v}\right|$, so that $2 \sigma_{0}$ satisfies the conditions stated in our theorem. Conversely, those conditions imply (i) and (ii) if $2 \sigma_{0}<2 n$.)

The case $2 \sigma_{0} \geq 2 n$ is similar; we use Theorem 17.12 (i) or (v) and (17.27); we need (ii) in order to insure that $\Psi(\mathbf{s}) \neq 0$; also, instead of (i) we assume
(iv) $\quad 2 n \leq 2 \sigma_{0} \leq m_{v}$ and $2 \sigma_{0}-\ell_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$.

Now $h_{v}+h_{v \rho}-\left|d_{v}\right|=m_{v}$ if $d_{v} \geq 0$, and $h_{v}+h_{v \rho}-\left|d_{v}\right|=m_{v}+2 d_{v}$ if $d_{v}<0$. Therefore, changing ( $h, m^{\prime}$ ) into ( $k, m$ ), we obtain the condition on $\sigma_{0}$ as stated in our theorem. (If $d_{v} \geq 0$, then (iv) implies that $2 h_{v \rho}+\ell_{v} \geq 2 n$, and hence $4 n-\left(2 h_{v \rho}+\ell_{v}\right) \leq 2 n \leq 2 \sigma_{0}$; if $d_{v}<0$, then $-\ell_{v}<h_{v \rho}-h_{v}$, and hence $4 n-\left(2 h_{v \rho}+\ell_{v}\right)<4 n-\left(h_{v}+h_{v \rho}\right)<n \leq 2 \sigma_{0}$, if we assume $h_{v}+h_{v \rho}>3 n$.) Also we have to exclude the case in which $2 \sigma_{0}=2 n+1, F=\mathbf{Q}$, and $\chi_{1}=\theta$. Even in that case we have (28.10) if $k_{v}-k_{v \rho} \neq \ell_{v}$, as will be shown by the second method.

Case SP (1st method). The case of integral $k$ can be proved in the same way as above. Writing again $h$ for the weight, take $e \in \mathbf{Z}^{\mathbf{a}}$ so that $\chi_{\mathbf{a}}(x)=x_{\mathbf{a}}^{h+e}\left|x_{\mathbf{a}}\right|^{-h-e}$ and $0 \leq e_{v} \leq 1$ for every $v \in \mathbf{a}$; put $k=h+e$. (That is what we did at the end of $\S 25.4$.) We evaluate (25.8a) at $s=\sigma_{0} / 2$. Suppose $\sigma_{0} \leq n$. We employ (17.22) and Theorem 17.7 (v) with $\nu=2 n+1-\sigma_{0}$. We have to assume that $0 \leq k_{v}-\nu \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$, that is, $2 n+1-h_{v}-e_{v} \leq \sigma_{0} \leq n$ and $\sigma_{0}-1+h_{v}-e_{v} \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$. In this case $\Psi(\mathbf{s}) \neq 0$ with no extra condition. Writing $(k, \mu)$ for ( $h, e$ ), we obtain "the left half" of our theorem with the condition as stated in our theorem. "The right half" can be obtained from (17.21) and Theorem 17.7 (iii). (If $2 n+1-k_{v}+\mu_{v} \leq \sigma_{0} \leq k_{v}-\mu_{v}$, then $k_{v}-\mu_{v} \geq n+1 / 2$ and $2 n+1-k_{v}+\mu_{v} \leq n+1 / 2$. Thus $\sigma_{0} \leq k_{v}-\mu_{v}$ if $\sigma_{0} \leq n$, and $\sigma_{0} \geq 2 n+1-k_{v}+\mu_{v}$ if $\sigma_{0}>n$.)

The case of half-integral weight can be handled by employing [S95b, (8.4)] instead of (25.8a).

Case SP (2nd method). Define $\mu \in \mathbf{Z}^{\mathbf{a}}$ and take $\sigma_{0}$ as in our theorem; assume that $\sigma_{0} \geq n+1 / 2$. Put $l=\mu+(n / 2) \mathbf{a}$ and $\nu=\sigma_{0}-(n / 2)$; then $\nu \geq(n+1) / 2$ and $0^{\bullet} \leq k-l-\nu \mathbf{a} \in 2 \mathbf{Z}^{\mathbf{a}}$. Define $\mathbf{g}$ as in §A5.5 with the present $\mu$ and $\chi$. Evaluating (28.4) at $s=(\nu-n-1) / 2$, we have again (28.8) with $Q=\Gamma(((\nu-n-1) / 2))^{-1} \Lambda\left(\sigma_{0}\right)$ for the present $\nu$ and $\sigma_{0}$. By the same procedure as in the proof of Theorem 28.5 we find an element $\mathbf{q}_{\nu}$ of $\mathcal{V}(\overline{\mathbf{Q}})$ such that

$$
\begin{equation*}
c_{\mathbf{f}}(\tau, r) \mathcal{Z}\left(\sigma_{0}\right)=\pi^{d_{0}+n|q|} \operatorname{det}(\tau)^{h+s \mathbf{a}} P_{\tau, r}\left(\sigma_{0}\right) Q\left\langle\mathbf{q}_{\nu}, \mathbf{f}\right\rangle \tag{28.11}
\end{equation*}
$$

with $q=(k-l-\nu \mathbf{a}) / 2$. We have $\operatorname{det}(\tau)^{h+s \mathbf{a}} \in \Phi \mathbf{Q}_{\mathrm{ab}}$ for the same reason as in the proof of Theorem 28.5. Let us now assume that $k \in \mathbf{Z}^{\mathbf{a}}$; then $\sigma_{0} \in \mathbf{Z}$. By (22.4a)
we have $q_{v}+s+h_{v}=k_{v}-(n+1) / 2$, and so from the explicit form of $\Gamma((s))$ in $\S 22.3$ we see that $\pi^{d_{0}+n|q|} \Gamma((s))^{-1} \in \pi^{n|k|-d \gamma} \mathbf{Q}^{\times}$, where $\gamma=n^{2} / 4$ if $n \in 2 \mathbf{Z}$ and $\gamma=\left(n^{2}-1\right) / 4$ if $n \notin 2 \mathbf{Z}$. Applying Lemma 17.5 (2) to each factor of $\Lambda$, we find that $\Lambda\left(\sigma_{0}\right) \in \pi^{d(n+1) \sigma_{0}-d \beta}$, where $\beta=n(3 n+4) / 4$ if $n \in 2 \mathbf{Z}$ and $\beta=(n+1)(3 n+4) / 4$ if $n \notin 2 \mathbf{Z}$. Therefore, dividing (28.11) by $\langle\mathbf{f}, \mathbf{f}\rangle$, we obtain our assertion in Case SP from Theorem 28.5, at least for $\sigma_{0} \geq n+1 / 2$. In view of Theorem 17.7 (i) we have to exclude the case $\nu=(n+2) / 2$ if $F=\mathbf{Q}$. Also, if $\nu=(n+3) / 2$, we cannot apply Lemma 15.8 to $E(\nu / 2)$. However, if $\mu \neq 0$, then $\mathbf{g}$ is a cusp form, and hence, by Theorem 17.9 (ia), $\pi^{-n|q|} e_{p} g_{p} E(\nu / 2)$ is a $\Phi \mathbf{Q}_{\mathrm{ab}}$-rational element of $\mathcal{R}_{\omega}^{a}$ of (15.4) with suitable $\omega$ and $a$, so that we can take $h_{p}=\pi^{-n|q|} \mathfrak{p}\left(e_{p} g_{p} E(\nu / 2)\right)$ with the map $\mathfrak{p}$ of Proposition 15.6 (3), and eventually obtain $\mathbf{q}_{\nu}$ from $h_{p}$. This is why we have the condition $[k]-\ell \in 2 \mathbf{Z}$ in the bad case (B). The case of half-integral weight can be handled in the same way.

Case UT (2nd method). Define $\mu \in \mathbf{Z}^{\mathbf{b}}$ by (22.15a, b) with $t^{\prime}=-\ell$; put $m^{\prime}=$ $\left(\mu_{v}+\mu_{v \rho}+n\right)_{v \in \mathbf{a}}$ and $\nu=2 \sigma_{0}-n$. Then $n \leq \nu \leq m_{v}-m_{v}^{\prime}$ and $m_{v}-m_{v}^{\prime}-\nu \in 2 \mathbf{Z}$ for every $v \in \mathbf{a}$. We have (28.4) with $\mathbf{g}$ defined by using the present $\mu$ and $\chi$. We evaluate (28.4) at $s=(\nu / 2)-n$; then we have $\operatorname{det}(\tau)^{h+s a} \in \Phi$ again. We see that $\chi_{\mathbf{a}}(x)=\operatorname{sgn}(x)^{2 \sigma_{0} \mathbf{a}}$ for $x \in F_{\mathbf{a}}^{\times}$, so that we can apply Lemma 17.5 (2) to $\Lambda\left(\sigma_{0}\right)$. Therefore we obtain our assertion in Case UT, at least for $\sigma_{0} \geq n$, in the same manner as in Case SP. There is no problem if $\mathbf{g}$ is a cusp form, and so the conditions in Case (E) are stated as above.

We now consider the arithmeticity of the Eisenstein series of Section 23.
28.9. Theorem (Cases SP and UT). Let $n, r, \mathbf{f}, g, \chi$ be as in Theorem 23.11; let $m_{0}$ be as in Theorem 28.5. Suppose that $\mathbf{f} \in \mathcal{V}(\overline{\mathbf{Q}})$ and $g$ is $\overline{\mathbf{Q}}$-rational; suppose also that $m_{0}>(3 r / 2)+1$ in Case $S P$ and $m_{0}>3 r$ in Case UT. (Notice that $\mathbf{f} \in \mathcal{S}_{k}^{r}(C)$, so that $n$ in Theorem 28.5 is now r.) Let $\lambda_{n}=(n+1) / 2$ in Case $S P$ and $\lambda_{n}=n$ in Case UT. In the setting of Theorem 23.11 we can state the arithmeticity of each function as follows:
(I) Let $\mu$ be as in Theorem 23.11 (I). Suppose that $\mu>(3 r / 2)+1$ in Case SP and $\mu>3 r$ in Case UT. If $F=\mathbf{Q}$, suppose moreover that $\mu \notin\left\{\lambda_{n+r}+1 / 2, \lambda_{n+r}+\right.$ 1\} in Case $S P$ and $\mu \neq \lambda_{n+r}+1$ in Case UT. Then $\pi^{-\alpha} E_{k}^{n, r}(z, \mu / 2 ; g, \Gamma)$ and $\pi^{-\alpha} E_{q}(z, \mu / 2 ; \mathbf{f}, \chi, C)$ are $\overline{\mathbf{Q}}$-rational, where $\alpha=\sum_{v \in \mathbf{a}}\left(m_{v}-\mu\right)(n-r) / 2$.
(II) Suppose $F=\mathbf{Q}$ and $\mu=\lambda_{n+r}+1$. Then $\pi^{-\alpha} E_{q}(z, \mu / 2 ; \mathbf{f}, \chi, C)$ with $\alpha$ as above is $\overline{\mathbf{Q}}$-rational in Case SP if $\chi^{2} \neq 1$, and also in Case UT if $\chi_{1} \neq \theta^{\mu}$.
(III) Let $\mu$ be as in Theorem 23.11 (III). Exclude Cases (A), (B), (C), (D), and $(F)$ of the same theorem. Then $\pi^{-\beta}\langle\mathbf{f}, \mathbf{f}\rangle^{-1} \mathcal{F}_{q}(z, \mu / 2 ; \mathbf{f}, \chi, C)$ is $\overline{\mathbf{Q}}$-rational, where $\beta=\sum_{v \in \mathbf{a}}\left(m_{v}+\mu\right)(n+r) / 2-d e$ with $e=\left[(n+r)^{2} / 4\right]$ in Case SP and $e=(n+r)(n+r-1) / 2$ in Case UT.

Proof. We evaluate (24.29) at $s=\mu / 2$ as we did in §25.7. By (25.11), $c_{m}(\mathbf{s}) \in$ $\pi^{d_{0}} \mathbf{Q}^{\times}$at $s=\mu / 2$, and by Lemma $17.5(2), \Lambda_{\mathrm{c}}^{2 r}(\mu / 2, \chi) \in \pi^{d M} \mathbf{Q}_{\mathrm{ab}}$ with an integer $M$ which can be explicitly given. Therefore, by (25.9) and (25.10), we find that

$$
\begin{equation*}
\mathcal{Z}(u \mu / 2, \mathbf{f}, \chi) E_{q}(z, \mu / 2 ; \mathbf{f}, \chi, C)=\pi^{\gamma} \sum_{a \in \mathcal{B}}\left\langle h_{a i}^{\prime}, f_{a}\right\rangle g_{a i}^{\prime}(z), \tag{28.12}
\end{equation*}
$$

where $\gamma=\sum_{v \in \mathbf{a}}\left(m_{v}-\mu\right)(n+r) / 2+d M$ and $g_{a i}^{\prime}, h_{a i}^{\prime} \in \mathcal{N}_{k}^{t}(\Psi)$ with some $t$. Now (25.9) was obtained by applying Lemma 24.11 to $\left(H_{q . a}\right)^{\circ}$. Instead, by virtue of Lemma 28.7, we can replace $J_{q . a}$ by a function $T_{q . a}(z, w)$, which is holomorphic
in $w$, and similar to $T_{b, a}$ in the proof of Theorem 28.8. We eventually find (28.12) with $h_{a i}^{\prime} \in \mathcal{M}_{k}(\Psi)$. Here we have to verify that $J_{q, a}$ is a function of type $\Delta_{\nu \mathrm{a}}^{p} R$ of Lemma 28.7, which is so if we exclude the bad cases stated in (I) and (II). Then we find that $\left\langle h_{a i}^{\prime}, f_{a}\right\rangle \in\langle\mathbf{f}, \mathbf{f}\rangle \overline{\mathbf{Q}}$ by Corollary 28.6. By Theorem 20.13 and (28.10), $\mathcal{Z}\left(\sigma_{0}, \mathbf{f}, \chi\right) \in \pi^{r|m|+d \varepsilon}\langle\mathbf{f}, \mathbf{f}\rangle \overline{\mathbf{Q}}^{\times}$for $\sigma_{0}=u \mu / 2$ with $\varepsilon$ given by the formula there in which we have to take $r$ in place of $n$. An explicit calculation of $M$ shows that $M=r \mu+\varepsilon$. Therefore we obtain our assertions for $E_{q}$ in (I) and (II). We obtain the assertion for $E_{k}^{n, r}$ for the same reason as in §25.7.

As for (III), the desired result can be derived from (24.29a), Theorem 17.7 (iii, $\mathrm{v})$, Theorem 17.12 ( $\mathrm{i}, \mathrm{iii}$ ), (17.22), and (17.30) in a similar way. If $k \notin \mathbf{Z}^{\mathbf{a}}$, then we use [S95b, (7.22)].
28.10. In the above theorem the case in which $F=\mathbf{Q}$ and $\mu=\lambda_{n+r}+1$ is excluded, though that case can be handled under a certain condition on $\chi$. It is conjecturable that the arithmeticity as in the above theorem is always true for $\mu=\lambda_{n+r}+1$ even when $F=\mathbf{Q}$. In fact, we can at least prove the rationality as stated in the above theorem if the following inequality holds:

$$
\begin{equation*}
(n+r)\left(\lambda_{n+r}-1\right)>4 \lambda_{r}-2+(n+r-2) m \tag{28.13}
\end{equation*}
$$

Here we assume $F=\mathbf{Q}$, and so $m=\left(m_{v}\right)_{v \in \mathbf{a}} \in 2^{-1} \mathbf{Z}$. Indeed, if $\mu=\lambda_{n+r}+1$, then the degree of near holomorphy, written $t$ in the above proof, is given by $t=(n+r)(m-\mu+2) / 2$. Then Lemma 28.2 is applicable to $h_{a i}^{\prime}$ of (28.12) if $m>2 \lambda_{r}-1+t$, which is equivalent to (28.13). Thus, under (28.13), that lemma allows us to replace $h_{a i}^{\prime}$ of (28.12) with $\mathfrak{p}\left(h_{a i}^{\prime}\right)$, which belongs to $\mathcal{S}_{k}^{r}(\Psi)$. Therefore we obtain the desired arithmeticity.
28.11. Corollary. The notation being as in Theorem 28.9, suppose that $m=$ $\mu \mathbf{a}$ with $\mu \in 2^{-1} \mathbf{Z}$ such that $\mu>(3 r / 2)+1$ in Case SP, $\mu>3 r$ in Case UT, and $\mu \geq \lambda_{n+r}$ in both cases; if $F=\mathbf{Q}$, suppose in addition that $\mu \notin\left\{\lambda_{n+r}+\right.$ $\left.1 / 2, \lambda_{n+r}+1\right\}$ in Case SP, and $\mu \neq \lambda_{n+r}+1$ in Case UT. Then $E_{k}^{n, r}(z, \mu / 2 ; g, \Gamma)$ and $E_{q}(z, \mu / 2 ; \mathbf{f}, \chi, C)$ are $\overline{\mathbf{Q}}$-rational.

Proof. This follows immediately from Theorem 28.9 (I), since $\alpha=0$ if $m=\mu \mathrm{a}$.
28.12. Proof of Theorem 27.16 (1) in Cases $S P$ and UT. Let us first treat Case SP. The desired fact is included in Theorem 28.9, but the proof of the theorem requires the map $\mathfrak{q}$ of Theorem 27.14, which is guaranteed by Theorem 27.16. Therefore, to prove Theorem 27.16, we need the same theorem, but this does not produce any vicious circle, since the proof of Theorem 28.9 on $G^{n}$ requires Theorem 27.16 on $G^{r}$ with $r<n$, and we can justify the whole proof by induction on the dimensionality, which is as follows.

First of all, the case $r=0$ of Theorem 27.16 (1) is guaranteed by Theorem 17.7 (i). Therefore, if $n=1$, the stability of $\mathcal{E}_{k}^{\psi}$ under $\operatorname{Aut}(\mathbf{C} / \overline{\mathbf{Q}})$ is true at least for $\mu>2$; in fact, it is true even for $\mu=2$ as noted in Remark 27.17 (A). Thus we have the map $\mathfrak{q}$ of Theorem 27.14 for the forms on $G^{1}$; see [S87b, Theorem 9.1 and Proposition 9.4] for the most comprehensive results in the case $n=1$. Let $\mu$ be as in Theorem 27.15 with $1=r<n$. Notice that $\mu \in \Lambda(1, \mu \mathbf{a})$ if and only if $\mu \geq 2$. Therefore $\mathcal{Z}(s, \mathbf{f}, \chi) \neq 0$ at $s=\mu$ if $\mu \geq 2$, and hence the proof of Theorem 28.5, as well as that of Corollary 28.6, is valid for the forms of weight $\mu \mathbf{a}$ on $G^{1}$ with such a $\mu$. Consequently (28.10) is valid for $n=1$ and $\sigma_{0}=\mu$; then our proof
of Theorem 28.9 establishes Theorem 27.16 (1) for $r=1$. This settles Theorem 27.16 for $n=2$. To prove Theorem 27.16 for an arbitrary $n>2$, we assume that it is true for $G^{n^{\prime}}$ for every $n^{\prime}<n$. Let $1<r<n$ and let $\mu$ be as in Theorem 27.15. Our induction assumption guarantees the map $\mathfrak{q}$ for the forms on $G^{r}$. Since $\mu>(3 r / 2)+1$, Theorem 28.5 and Corollary 28.6 are valid for the forms of weight $\mu \mathrm{a}$ on $G^{r}$, and again the proof of Theorem 28.9 establishes Theorem 27.16 (1) for those $r$ and $\mu$. This completes the proof in Case SP. Case UT can be handled similarly.
28.13. Remark. (A) For $m=\mu \mathbf{a}$ we already stated in Theorem 27.16 the arithmeticity of $E_{k}^{n, r}(z, \mu / 2 ; g, \Gamma)$ under the condition that $\mu>\lambda_{n+r}$ if $F \neq \mathbf{Q}$ and $\mu>\lambda_{n+r}+1$ if $F=\mathbf{Q}$. Thus the result for $\mu=\lambda_{n+r}$ given in Corollary 28.11 is not included in that theorem.
(B) In [S76] and several subsequent papers the author obtained various results concerning the critical values of the zeta functions on $G L_{2}$ and $G L_{2} \times G L_{2}$, as well as some other related zeta functions. The most comprehensive results in the case with a totally real number field as the basic field were given in [S88] and [S91], in which references to the papers in the intermediate period can be found. In [S76] we formulated the result of type (28.10) in the form, for example,

$$
\begin{equation*}
\left[\pi^{-k}\langle f, f\rangle^{-1} D(m, f, g)\right]^{\sigma}=\pi^{-k}\left\langle f^{\sigma}, f^{\sigma}\right\rangle^{-1} D\left(m, f^{\sigma}, g^{\sigma}\right) \tag{28.14}
\end{equation*}
$$

for every $\sigma \in \operatorname{Aut}(\mathbf{C})$, where $f$ is a Hecke eigenform of weight $k$ and $g$ is another modular form with respect to congruence subgroups of $S L_{2}(\mathbf{Z}) ; D(s, f, g)$ is of type (22.4). Though we can in fact state (28.10) in such a form, in the present book we content ourselves with a weaker statement, since the proof of the results in the form (28.14) would make our exposition longer and more tedious. At any rate we believe that what was done in Section 10 combined with careful examinations of the behavior of the Eisenstein series of Sections 16 and 17 under Aut(C) can give the desired formulas with no extra new ideas. Also, we can prove, employing (22.9), the algebraicity of the critical values of the functions of (22.4) in the form similar to (28.10), but the task of giving precise statements for these can be left to the reader as easy exercises.

There are a few more technical points. In [S76] we relied on the existence of a "primitive" Hecke eigenform, which is not guaranteed in general, in the case of $S p(n, F)$, for example. However, in [S81a] we introduced a method by which this difficulty can be avoided, and mentioned that the higher-dimensional symplectic case could be handled by the same method. Indeed, in this section we proved the expected results by the basic idea of [S76] combined with the technique of [S81a]. The latter requires the nonvanishing of $\mathcal{Z}$ at a critical point, which we proved in Theorem 20.13, and the result is best possible in the sense of Theorem 22.11. However, given a point $\sigma$ smaller than the limit of Theorem 20.13, it seems possible to find a suitable $\chi$ such that $\mathcal{Z}(\sigma, \mathbf{f}, \chi) \neq 0$. A result of this nature for the zeta function on $G L_{2}$ was proved by Rohrlich in $[\mathrm{R}]$. A similar result in the higher-dimensional case will certainly allow us to state our theorems in improved forms.
(C) As we already noted at the end of Section 15 , if $G=S L_{2}(F)$, the map $\mathfrak{p}$ of Lemma 28.2 can be established for every $k \in 2^{-1} \mathbf{Z}^{\mathbf{a}}$ with no condition on the degree of near holomorphy, and the nonvanishing of $\mathcal{Z}$ can be given in a better form, as noted in Theorem 22.13. Therefore Theorems 28.5 and 28.8 can be stated
in stronger forms. For example, if $k \in \mathbf{Z}^{\mathbf{a}}$, these theorems are true if we assume $m_{0} \geq 2$ instead of (28.7). For details, see [S91] and the papers cited there.

Also, take $r=1$ in the setting of Theorem 28.9 in Case SP and assume $m_{0} \geq 2$. Then the conclusions of (I) and (II) are true for every $\mu$ as in Theorem 23.11 (I) and (III), excluding the cases specified there, as already noted in $\S 28.12$.

## 29. Main theorems on arithmeticity in Case UB

29.1. We now consider Case UB with $(V, \varphi), G^{\varphi}$, and $\mathcal{Z}^{\varphi}$ as in $\S 26.1$; we use the symbols $r_{v}, t_{v}$, and $\mathbf{i}^{\varphi}$ defined by (26.1) and (26.5). Also, we put $d=[F$ : $\mathbf{Q}], n=\operatorname{dim}(V)$, and

$$
\begin{equation*}
\mathbf{a}^{\prime}=\left\{v \in \mathbf{a} \mid r_{v}>0\right\} \tag{29.1}
\end{equation*}
$$

This is the same as $\mathbf{a}^{\prime}$ of (12.5). By (3.23) and (3.24a, b) we have

$$
\begin{equation*}
j_{\alpha}^{k}(z)=\prod_{v \in \mathbf{a}} \operatorname{det}\left(\alpha_{v}\right)^{-k_{v \rho}} \prod_{v \in \mathbf{a}^{\prime}} j_{v}(\alpha, z)^{k_{v}+k_{v \rho}} \quad\left(k \in \mathbf{Z}^{\mathbf{b}}, \alpha \in G_{\mathbf{A}}^{\varphi}\right) \tag{29.2}
\end{equation*}
$$

Thus we can ignore the $k_{v}$ for $v \in \mathbf{a}^{\prime}$ in the definition of $\mathcal{M}_{k}^{\varphi}$, but we need them in the definition of Eisenstein series of type $E_{k}^{\psi, \varphi}$.

In our proof it will become necessary to study the nature of $\left(A_{m}^{k} f\right)(z, w)$ defined by [S97, (23.6.4)]. To recall the definition of $A_{m}^{k} f$, we take $\psi=\varphi$ in the setting of $\S 26.11$. Then both $\iota$ and $\iota_{U}$ are maps of $\mathfrak{Z}^{\varphi} \times \mathcal{Z}^{\varphi}$ into $\mathcal{H}_{n}^{\mathbf{a}}$. Let $k \in \mathbf{Z}^{\mathbf{b}}$ and $q \in \mathbf{Z}^{\mathbf{a}^{\prime}}$. (This $q$ is not $q$ in $\S 26.11 ; q$ there is now 0 ; also we use $q$ in place of $m$ in $[\mathrm{S} 97,(23.6 .4)]$, so that we shall speak of $A_{q}^{k}$.) We consider $S_{p}(T)$ of $\S 13.13$ for Type A with $p \in \mathbf{Z}^{\mathbf{a}^{\prime}}$ and $T=\prod_{v \in \mathbf{a}^{\prime}} T_{v}, T_{v}=\mathbf{C}_{n}^{n}$ for every $v \in \mathbf{a}^{\prime}$. Then, for $u \in T$ and $(z, w) \in \mathfrak{Z}^{\varphi} \times \mathfrak{Z}^{\varphi}$ we put

$$
\begin{equation*}
\xi(u)=\prod_{\substack{v \in \mathbf{a} \\ q_{v}>0}} \operatorname{det}\left[\ell_{v}\left(u_{v}\right)\right]^{q_{v}}, \quad \zeta_{z, w}(u)=\prod_{\substack{v \in \mathbf{a} \\ q_{v}<0}} \operatorname{det}\left[\ell_{v}^{\prime}\left({ }^{t} M\left(w_{v}\right) u_{v} N\left(z_{v}\right)\right)\right]^{\left|q_{v}\right|} \tag{29.3}
\end{equation*}
$$

Here $\ell_{v}(Y)$ (resp. $\ell_{v}^{\prime}(Y)$ ) denote the lower left (resp. the lower right) ( $r_{v} \times r_{v}$ )block of $Y ; M\left(w_{v}\right)$ and $N\left(z_{v}\right)$ are defined by [S97, (6.11.4)]. Then $\xi \in S_{p^{\prime}}(T)$ and $\zeta_{z, w} \in S_{p^{\prime \prime}}(T)$ with $p^{\prime}=\left(\operatorname{Max}\left\{r_{v} q_{v}, 0\right\}\right)_{v \in \mathbf{a}^{\prime}}$ and $p^{\prime \prime}=\left(\operatorname{Max}\left\{-r_{v} q_{v}, 0\right\}\right)_{v \in \mathbf{a}^{\prime}}$; there is an irreducible subspace $W$ (resp. $Z$ ) of $S_{p^{\prime}}(T)$ (resp. $S_{p^{\prime \prime}}(T)$ ) containing $\xi\left(\operatorname{resp} \zeta_{z, w}\right)$; see [S97, §23.6]. Now for a function $f$ on $\mathcal{H}_{n}^{\mathrm{a}}$ we define $A_{q}^{k} f$ to be a function on $\mathfrak{Z}^{\varphi} \times \mathfrak{Z}^{\varphi}$ given by

$$
\begin{equation*}
\left(A_{q}^{k} f\right)(z, w)=B_{z, w}\left(f \|_{k} U^{-1}\right)(\iota(z, w)), \quad B_{z, w} g=\left(E^{Z} D_{\rho}^{W} g\right)\left(\xi, \zeta_{z, w}\right) \tag{29.4}
\end{equation*}
$$

where $U$ is as in [S97, (22.1.6)], $g$ is a function on $\mathcal{H}_{n}^{\mathbf{a}}, \rho(x)=\operatorname{det}(x)^{k}$, and $E^{Z} D_{\rho}^{W} g$, as defined in $\S 13.13$, is a function on $\mathcal{H}_{n}^{\mathrm{a}}$ with values in $S_{1}\left(W, S_{1}(Z, \mathbf{C})\right)$. In [S97, Lemma 23.9] we showed

$$
\begin{equation*}
\prod_{v \in \mathbf{a}} \operatorname{det}(\gamma)^{k_{v}-k_{v \rho}} A_{q}^{k}\left(f \|_{k}[\beta, \gamma]_{S}\right)(z, w)=j_{\beta}^{k+q}(z)^{-1}{\overline{j_{\gamma}^{k+q}(w)}}^{-1}\left(A_{q}^{k} f\right)(\beta z, \gamma w) \tag{29.5}
\end{equation*}
$$

for $(\beta, \gamma) \in G^{\varphi} \times G^{\varphi}$, where $[\beta, \gamma]_{S}$ is defined by (26.36). (The reader is reminded of the difference between the present notation about the factor of automorphy and that of [S97]; see §5.4.)
29.2. Lemma. Let $f \in \mathcal{N}_{k}^{e}(\overline{\mathbf{Q}})$ with $e \in \mathbf{Z}^{\mathbf{a}}$; put $h=k+q, \alpha=\sum_{v \in \mathbf{a}} r_{v} q_{v}$, and $e^{\prime}=\left(e_{v}^{\prime}\right)_{v \in \mathbf{a}}, e_{v}^{\prime}=\operatorname{Max}\left\{e_{v}+r_{v} q_{v}, 0\right\}$. Then

$$
\begin{equation*}
\pi^{-\alpha} \mathfrak{q}^{-1}\left(A_{q}^{k} f\right)(z, w)=\sum_{i \in I} g_{i}(z) \overline{h_{i}(w)} \tag{29.6}
\end{equation*}
$$

with a finite set of indices $I$ and $g_{i}, h_{i} \in \mathcal{N}_{h}^{e^{\prime}}(\overline{\mathbf{Q}})$, where $\mathfrak{q}=p_{K}\left(\sum_{v \in \mathbf{a}}\left(k_{v}-\right.\right.$ $\left.\left.k_{v \rho}\right) \tau_{v}, \sum_{v \in \mathbf{a}} t_{v} \tau_{v}\right)$.

Proof. By (29.5), $A_{q}^{k} f$ has at least the desired automorphy property of the right-hand side of (29.6). To prove its near holomorphy, we can reduce the problem to the nature of $\left(A_{q}^{k} f\right)\left(z, \mathbf{i}^{\varphi}\right)$ or $\left(A_{q}^{k} f\right)\left(\mathbf{i}^{\varphi}, w\right)$ by the same argument as in the proof of Lemma 26.12. Put $g=f \| U^{-1}$. Then $g \in \mathcal{N}^{e}$, and so, by Proposition 13.15 (1), the components of $E^{Z} D_{\rho}^{W} g$ belong to $\mathcal{N}^{e^{\prime}}$. Since $\zeta_{z, w}$ and $\iota(z, w)$ are holomorphic in $(z, \bar{w})$, we see that $\left(A_{q}^{k} f\right)(z, w)$ is a polynomial in $\left(\iota(z, w)-\iota(z, w)^{*}\right)^{-1}$ whose coefficients are holomorphic functions in $(z, \bar{w})$. Therefore the argument in the proof of Lemma 26.12 shows the desired near holomorphy.

To prove the $\overline{\mathbf{Q}}$-rationality, take the symbols $z_{0}, w_{0}, \mathfrak{z}_{0}$, and $\mathfrak{z}_{1}$ as in $\S 26.11$. Put $\lambda=\lambda\left(U^{-1}, \mathfrak{z}_{1}\right), \mu=\mu\left(U^{-1}, \mathfrak{z}_{1}\right)$, and $j=j\left(U^{-1}, \mathfrak{z}_{1}\right)$ for simplicity. From the definition of $U$ in [S97, (22.1.6)] we see that $U_{v}$ has algebraic entries for every $v \in \mathbf{a}$; also by Lemma $4.13, z_{0}, w_{0}, z_{1}$, and $z_{0}$ have algebraic coordinates. Therefore $\xi, \zeta_{z_{0}, w_{0}}, M\left(w_{0}\right), N\left(z_{0}\right), \lambda, \mu$, and $j$ are all $\overline{\mathbf{Q}}$-rational. Define $\xi^{\prime}$ and $\zeta^{\prime}$ by $\xi^{\prime}(u)=\xi\left(\lambda^{-1} u \cdot{ }^{t} \mu^{-1}\right)$ and $\left.\zeta^{\prime}(u)=\zeta_{z_{0}, w_{0}}{ }^{t} \lambda u \mu\right)$. Then $\left(B_{z_{0}, w_{0}} g\right)\left(\mathfrak{z}_{1}\right)=$ $\left(E^{Z} D_{\rho}^{W} g\right)\left(\xi, \zeta_{z_{0}, w_{0}}\right)\left(\mathfrak{z}_{1}\right)=j^{-k}\left(E^{Z} D_{\rho}^{W} f\right)\left(\xi^{\prime}, \zeta^{\prime}\right)\left(\mathfrak{z}_{0}\right)$ by the generalization of (12.21) and (12.24a, b) mentioned in §13.13. Put $\xi_{1}(u)=\xi^{\prime}\left(\mathfrak{p}_{v \rho}\left(\mathfrak{z}_{0}\right)^{-1} u \cdot{ }^{t} \mathfrak{p}_{v}\left(\mathfrak{z}_{0}\right)^{-1}\right)$ and $\zeta_{1}(u)=\zeta^{\prime}\left({ }^{t} \mathfrak{p}_{v \rho}\left(\mathfrak{z}_{0}\right) u \mathfrak{p}_{v}\left(\mathfrak{z}_{0}\right)\right)$. Fixing our attention on one $v \in \mathbf{a}$, we suppress the subscript $v$. By (26.41a, b)) we have

$$
\begin{aligned}
\zeta_{1}(u) & =\operatorname{det}\left[\ell^{\prime}\left({ }^{t} M\left(w_{0}\right) \cdot{ }^{t} \lambda \cdot{ }^{t} \mathfrak{p}_{v \rho}\left(\mathfrak{z}_{0}\right) u \mathfrak{p}_{v}\left(\mathfrak{z}_{0}\right) \mu N\left(z_{0}\right)\right)\right]^{-q} \\
& =\operatorname{det}\left[\ell^{\prime}\left(\operatorname{diag}\left[t^{t} \mathfrak{p}_{\rho}\left(z_{0}\right),{ }^{t} \mathfrak{p}\left(w_{0}\right)\right] \cdot{ }^{t} M\left(w_{0}\right) \cdot{ }^{t} \lambda u \mu N\left(z_{0}\right) \operatorname{diag}\left[\mathfrak{p}_{\rho}\left(w_{0}\right), \mathfrak{p}\left(z_{0}\right)\right]\right)\right]^{-q} \\
& =\operatorname{det}\left(\mathfrak{p}\left(w_{0}\right) \mathfrak{p}\left(z_{0}\right)\right)^{-q} \zeta^{\prime}(u) .
\end{aligned}
$$

From the explicit forms of $M(w)$ and $N(z)$ in [S97, (6.11.4)] we see that $\ell\left(M\left(w_{0}\right) u\right.$. $\left.{ }^{t} N\left(z_{0}\right)\right)=\ell^{\prime}(u)$. Therefore we obtain $\xi_{1}=\operatorname{det}\left(\mathfrak{p}\left(w_{0}\right) \mathfrak{p}\left(z_{0}\right)\right)^{-q} \xi^{\prime}$ by the same type of calculation as for $\zeta_{1}$ and $\zeta^{\prime}$. Now, in the proof of Lemma 26.12 we have shown that $\mathfrak{P}_{k}\left(\mathfrak{z}_{0}\right)=\mathfrak{P}_{k}\left(z_{0}\right) \mathfrak{P}_{k}\left(w_{0}\right) \mathfrak{q}$. Therefore

$$
\begin{aligned}
\left(\mathfrak{P}_{\rho \otimes \tau^{W} \otimes \sigma^{z}}\left(\mathfrak{z}_{0}\right)^{-1} E^{Z} D_{\rho}^{W} f\right)\left(\xi^{\prime}, \zeta^{\prime}\right)\left(\mathfrak{z}_{0}\right) & =\left[\mathfrak{P}_{k}\left(z_{0}\right) \mathfrak{P}_{k}\left(w_{0}\right) \mathfrak{q}\right]^{-1}\left(E^{Z} D_{\rho}^{W} f\right)\left(\xi_{1}, \zeta_{1}\right)\left(\mathfrak{z}_{0}\right) \\
& =\left[\mathfrak{P}_{h}\left(z_{0}\right) \mathfrak{P}_{h}\left(w_{0}\right) \mathfrak{q}\right]^{-1}\left(E^{Z} D_{\rho}^{W} f\right)\left(\xi^{\prime}, \zeta^{\prime}\right)\left(\mathfrak{z}_{0}\right) \\
& =j^{k}\left[\mathfrak{P}_{h}\left(z_{0}\right) \mathfrak{P}_{h}\left(w_{0}\right) \mathfrak{q}\right]^{-1}\left(B_{z_{0}, w_{0}} g\right)\left(\mathfrak{z}_{1}\right)
\end{aligned}
$$

By Theorem $14.9(2), \pi^{-\alpha} E^{Z} D_{\rho}^{W} f \in \mathcal{N}_{\rho \otimes \tau^{W}{ }_{\otimes \sigma^{Z}}}^{e^{\prime}}(\overline{\mathbf{Q}})$. Since $\xi^{\prime}$ and $\zeta^{\prime}$ are $\overline{\mathbf{Q}}$ rational, $\pi^{-\alpha}$ times the first quantity of the above series of equalities is algebraic, and hence $\pi^{-\alpha}\left[\mathfrak{P}_{h}\left(z_{0}\right) \mathfrak{P}_{h}\left(w_{0}\right)\right]^{-1} \mathfrak{q}^{-1}\left(A_{q}^{k} f\right)\left(z_{0}, w_{0}\right) \in \overline{\mathbf{Q}}$. Now we can repeat the proof of Lemma 26.12 with $A_{q}^{k} f$ in place of $f^{\circ}$ there, since the necessary properties of the function are guaranteed by what we proved in the above. Thus we obtain (29.6).
29.3. Lemma. In the setting of $\S 26.11$, let $R \in \mathcal{M}_{\nu \mathbf{a}}^{\eta}(\overline{\mathbf{Q}}), S_{0}=B\left(\left(\Delta_{\nu \mathbf{a}}^{p} R\right) \|_{k}\right.$ $\left.U^{-1}\right)$, and $S(z, w)=S_{0}(\iota(z, w))$ for $(z, w) \in \mathfrak{Z}^{\psi} \times \mathfrak{Z}^{\varphi}$, where $\nu \in \mathbf{Z}, 0 \leq p \in \mathbf{Z}^{\mathbf{a}}$, $k$ is an element of $\mathbf{Z}^{\mathbf{b}}$ such that $\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}=2 p+\nu \mathbf{a}, \Delta_{\nu \mathbf{a}}^{p}$ is as in Lemma 15.8 and the proof of Theorem 17.9 in the unitary case, and $B$ denotes the operator $B_{z, w}$ of (29.4) with $q_{v} \geq 0$ for every $v$, which we consider only if $\psi=\varphi$. (Thus $B g=$ $\left(D_{\rho}^{W} g\right)(\xi)$, which does not involve the parameters $(z, w)$, and $S_{0}=\Delta_{\nu \mathbf{a}}^{p}\left(R \|_{\nu \mathbf{a}} U^{-1}\right)$ if $\psi \neq \varphi$.) Let $h=k+e-e^{\prime}$ and $m=\left(h_{v}+h_{v \rho}\right)_{v \in \mathbf{a}}$; let $L_{v}$ be the operator of (15.3) defined on $\mathfrak{Z}^{\varphi}$ with $\omega(a, b)=\operatorname{det}(b)^{m}$. Then there exists an element $T(z, w)$ of $\sum_{v \in \mathbf{a}} \sum_{i=0}^{\infty} \overline{\mathbf{Q}} L_{v}^{i} \bar{S}$ which is holomorphic in $w$ and such that $\langle\overline{S(z, w)}, f(w)\rangle=$ $\langle T(z, w), f(w)\rangle$ for every $f \in \mathcal{S}_{h}^{\varphi}$, provided $\nu \geq \operatorname{Max}_{v \in \mathbf{a}} r_{v}$.

This will be proven in $\S$ A8.13. Once this is established, we have an expression

$$
\begin{equation*}
\pi^{-N|p|-\alpha} \mathfrak{q}^{-1} T(z, w)=\sum_{i \in I} \overline{g_{i}(z)} h_{i}(w) \tag{29.7}
\end{equation*}
$$

where $\alpha$ and $\mathfrak{q}$ are as in (29.6), $I$ is a finite set of indices, $N=2^{-1}(\operatorname{dim}(V)+$ $\operatorname{dim}(W)), g_{i} \in \mathcal{N}_{h}^{e^{\prime}}(\overline{\mathbf{Q}})$, and $h_{i} \in \mathcal{M}_{h}(\overline{\mathbf{Q}})$. Indeed, put $f=\pi^{-N|p|} \Delta_{\nu \mathbf{a}}^{p} R$. Then $f$ is $\overline{\mathbf{Q}}$-rational by Theorem 14.12, and $S=\pi^{N|p|} A_{q}^{k} f$, to which Lemma 29.2 is applicable. Thus $S$ has an expression of the form (29.6), so that $\bar{S}$ has an expression of the form (29.7) with $\overline{\mathbf{Q}}$-rational nearly holomorphic $g_{i}$ and $h_{i}$. By Theorem 14.9 (2), the same is true for every element of $\sum_{v \in \mathbf{a}} \sum_{i=0}^{\infty} \overline{\mathbf{Q}} I_{v}^{i} \bar{S}$, and hence, for $T$ in particular. Now $T$ is holomorphic in $w$, and so the proof of Lemma 26.12, modified in an obvious way, gives the desired expression (29.7).
29.4. The notation being as in $\S 26.10$, let $\{\lambda(\mathfrak{a})\}$ be a system of eigenvalues on $\mathcal{S}_{h}^{\varphi}\left(D^{\varphi}\right)$ in the sense that $\mathbf{f}_{0} \mid T(\mathfrak{a})=\lambda(\mathfrak{a}) \mathbf{f}_{0}$ for every $\mathfrak{a}$ with some $\mathbf{f}_{0} \in \mathcal{S}_{h}^{\varphi}\left(D^{\varphi}\right), \neq$ 0 , where $h \in \mathbf{Z}^{\mathbf{b}}$. Put

$$
\begin{array}{cc}
\mathcal{V}=\left\{\mathbf{f} \in \mathcal{S}_{h}^{\varphi}\left(D^{\varphi}\right)|\mathbf{f}| T(\mathfrak{a})=\lambda(\mathfrak{a}) \mathbf{f}\right. & \text { for every } \mathfrak{a}\} \\
\mathcal{V}(\overline{\mathbf{Q}})=\mathcal{V} \cap \mathcal{S}_{h}^{\varphi}\left(D^{\varphi}, \overline{\mathbf{Q}}\right) & \text { (see } \S 28.1) \tag{29.8b}
\end{array}
$$

By Lemma 26.14, the $\lambda(\mathfrak{a})$ are algebraic. Therefore we have $\mathcal{V}=\mathcal{V}(\overline{\mathbf{Q}}) \otimes_{\mathbf{Q}} \mathbf{C}$. For the same reason as in the proof of Theorem 28.5 we have $\mathcal{S}_{h}^{\varphi}\left(D^{\varphi}, \overline{\mathbf{Q}}\right)=\mathcal{V}(\overline{\mathbf{Q}}) \oplus \mathcal{U}$ with a vector space $\mathcal{U}$ over $\overline{\mathbf{Q}}$ that is orthogonal to $\mathcal{V}(\overline{\mathbf{Q}})$. We are going to state our main theorems on the arithmeticity, in which we need also a Hecke character $\chi$ of $K$ such that $\chi_{\mathbf{a}}(x)=x_{\mathbf{a}}^{\ell}\left|x_{\mathbf{a}}\right|^{-\ell}$ with $\ell \in \mathbf{Z}^{\mathbf{a}}$. Here $\ell$ is basically arbitrary; also we denote by $h$, instead of $k$, the weight of a Hecke eigenform. However, whenever we consider an Eisenstein series $E_{p}\left(z, \nu / 2 ; \mathbf{f}, \chi, D^{\psi}\right)$, we denote the weight by $k$ and assume that $\ell=\left(k_{v}-k_{v \rho}\right)_{v \in \mathbf{a}}$, as we did in §26.10.
29.5. Theorem. The notation being as above, put $m_{0}=\operatorname{Min}_{v \in \mathbf{a}^{\prime}}\left\{h_{v}+h_{v \rho}\right\}$ and assume that $m_{0}>2 n$ and $G_{a}^{\varphi}$ is not compact. (See Theorem 29.7 below for the result in the compact case.) Then

$$
\begin{equation*}
\left\langle\mathbf{g}, \mathbf{g}^{\prime}\right\rangle /\langle\mathbf{f}, \mathbf{f}\rangle \in \overline{\mathbf{Q}} \text { if } \mathbf{g}, \mathbf{g}^{\prime} \in \mathcal{V}(\overline{\mathbf{Q}}) \text { and } 0 \neq \mathbf{f} \in \mathcal{V}(\overline{\mathbf{Q}}) \tag{29.9}
\end{equation*}
$$

Moreover, let $0 \neq \mathbf{f} \in \mathcal{V}(\overline{\mathbf{Q}})$ and let $\sigma_{0}$ be an element of $2^{-1} \mathbf{Z}$ such that

$$
2 n-\left(2 h_{v \rho}+\ell_{v}\right) \leq 2 \sigma_{0} \leq \begin{cases}2 h_{v}-\ell_{v} & \text { if } r_{v}\left(h_{v}-h_{v \rho}-\ell_{v}\right)<0  \tag{29.10a}\\ 2 h_{v \rho}+\ell_{v} & \text { otherwise }\end{cases}
$$

(29.10c) $\quad \sigma_{0}<0$ or $2 \sigma_{0} \geq n$ if $\mathfrak{c}=\mathfrak{g}, \quad \chi_{1}=\theta^{2 \sigma_{0}}$, and the conductor of $\chi$ is $\mathfrak{r}$,
where $\chi_{1}$ is the restriction of $\chi$ to $F_{\mathbf{A}}^{\times}$and $\theta$ is the Hecke character of $K$ corresponding to $K / F$. If $\varphi$ is isotropic, suppose also that

$$
\begin{equation*}
h_{v}-h_{v \rho} \geq \ell_{v} \quad \text { for every } v \in \mathbf{a}, \tag{29.11a}
\end{equation*}
$$

Then

$$
\begin{equation*}
\mathcal{Z}\left(\sigma_{0}, \mathbf{f}, \chi\right) \in \pi^{\gamma} \mathfrak{q}\langle\mathbf{f}, \mathbf{f}\rangle \overline{\mathbf{Q}} \tag{29.12}
\end{equation*}
$$

where $\gamma=(d n / 2)\left(2 \sigma_{0}-n+1\right)+\sum_{v \in \mathbf{a}}\left\{r_{v}\left(h_{v}+h_{v \rho}\right)+\left(t_{v} / 2\right)\left(2 h_{v \rho}+\ell_{v}\right)\right\}$ and $\mathfrak{q}=p_{K}\left(\sum_{v \in \mathbf{a}} \ell_{v} \tau_{v}, \sum_{v \in \mathbf{a}} t_{v} \tau_{v}\right)$ with the CM-type $\tau=\sum_{v \in \mathbf{a}} \tau_{v}$ of $K$ fixed in §3.5 and the period symbol $p_{K}$ of $\S 11.3$.

Notice that $\gamma \in \mathbf{Z}$. Indeed, for every $v \in \mathbf{a}$ we have $n \sigma_{0}+\left(t_{v} / 2\right) \ell_{v}=2 r_{v} \sigma_{0}+$ $\left(t_{v} / 2\right)\left(2 \sigma_{0}+\ell_{v}\right) \in \mathbf{Z}$ by (29.10b), from which we can easily derive that $\gamma \in \mathbf{Z}$.

Proof. Put $q=\left(h_{v}-h_{v \rho}-\ell_{v}\right)_{v \in \mathbf{a}}$ and $k=h-q$; then $\ell=\left(k_{v}-k_{v \rho}\right)_{v \in \mathbf{a}}$. We needed the ideals $\mathfrak{b}$ and $\mathfrak{c}$ for the definition of $D^{\varphi}$. Changing $\mathfrak{c}$ for its suitable multiple, we may assume, without changing $\mathcal{Z}(s, \mathbf{f}, \chi)$, that the conductor of $\chi$ divides $\mathfrak{c}$. Define $E_{\mathbf{A}}$ by (16.27) in Case UT with these $k, \chi, \mathfrak{b}, \mathfrak{c}$. Then (16.24a, b) are satisfied with $\kappa=0$. In [S97, (23.11.3)] we proved

$$
\begin{align*}
& \varepsilon(s) c_{h}\left(\mathbf{s}^{\prime}\right) \Psi(\mathbf{s}) \mathcal{Z}(s, \mathbf{f}, \chi) f_{b}(z)  \tag{29.13}\\
& \quad=\sum_{a \in \mathcal{B}} \chi_{\mathbf{h}}(\operatorname{det}(a)) \int_{\mathcal{D}_{a}} \Lambda_{\mathfrak{c}}^{n}(s, \chi)\left(A_{q}^{k} H_{b, a}\right)(z, w) f_{a}(w) \delta(w)^{m} \mathbf{d} w .
\end{align*}
$$

Here $m=\left(h_{v}+h_{v \rho}\right)_{v \in \mathbf{a}}$; the symbols $\mathcal{B},\left(f_{b}\right)_{b \in \mathcal{B}}, \Lambda_{\mathfrak{c}}^{n}, H_{b, a}$, and $c_{h}$ are essentially the same as those in the proof of Theorem 26.13 , though we have $\psi=\varphi$ here; they are explained in [S97, $\S \S 23.10$ and 23.11]; $\Psi$ is given in [S97, Lemma 23.8]; $\varepsilon(s)$ is the factor written $e \cdot d^{s}$ in [S97, (23.11.3)], which equals $e(k, m, \mathbf{s}) c(\chi)$ of [S97, (23.10.3)]. (In fact, $k$ and $m$ there are $\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$ and $q$ here. Also we should add corrections: The symbol $\zeta\left(\mathfrak{y} \kappa_{\alpha}^{*} \cdot{ }^{t} \mu_{\alpha}^{-1}\right)$ of [S97, p.192, line 3 and line 9 from the bottom] should be $\zeta\left(-\mathfrak{y} \kappa_{\alpha}^{*} \cdot{ }^{t} \mu_{\alpha}^{-1}\right)$; also, the right-hand side of the equality in [S97, p.194, line 5 from the bottom] needs an extra factor $(-1)^{t}$, where $t$ is the sum of $r_{v} m_{v}$ for all $v \in \mathbf{a}$ such that $m_{v}>0$; the quantity of [S97, p.195, line 5] needs an extra factor ( -1$)^{r m}$.) Putting $m_{v}^{*}=k_{v}+k_{v \rho}$ and $\mu=2 \sigma_{0}$, we evaluate (29.13) at $s=\mu / 2$ with $\mu$ under the condition

$$
\begin{equation*}
2 n-m_{v}^{*} \leq \mu \leq m_{v}^{*} \text { and } \mu-m_{v}^{*} \in 2 \mathbf{Z} \text { for every } v \in \mathbf{a} . \tag{}
\end{equation*}
$$

Observe that $m_{v}^{*}=2 h_{v \rho}+\ell_{v}$ and $\left(^{*}\right)$ follows from (29.10a, b), since $m_{v}^{*}=2 h_{v}-$ $\ell_{v}-2 q_{v}$. Now from [S97, (A2.9.2)] we easily see that $c_{h}\left(\mathbf{s}^{\prime}\right) \in \pi^{d_{0}} \overline{\mathbf{Q}}^{\times}$at $s=\sigma_{0}$, where $d_{0}$ is the complex dimension of $\boldsymbol{Z}^{\varphi}$. As for $\Psi(\mathbf{s})$, by [S97, (23.2.1) and Lemma 23.8 ], its $v$-factor in the obvious sense equals

$$
\begin{equation*}
\prod_{i=1}^{r_{v}} \prod_{j=1}^{\left|q_{v}\right|}\left(-s+\left(m_{v}^{*} / 2\right)+i-j\right) \tag{}
\end{equation*}
$$

if $r_{v}>0$ and $q_{v}<0$. For such a $v$, we easily see that ( ${ }^{* *}$ ) is nonzero at $s=\sigma_{0}$ under (29.10a). Examining the case $r_{v} q_{v} \geq 0$ in a similar way, we find that $\Psi(\mathbf{s}) \in$ $\mathbf{Q}^{\times}$at $s=\sigma_{0}$ under (29.10a). Clearly $\varepsilon\left(\sigma_{0}\right) \in \overline{\mathbf{Q}}^{\times}$. Now $H_{b, a}$ is a function of type $E_{r}$ of (17.23a) obtained from $E_{\mathbf{A}}$, and hence $\Lambda_{\mathfrak{c}}^{n}(s, \chi) H_{b, a}$ is a function of type $D_{r}$ of (17.24). Therefore, by Theorem 17.12 (v), its value at $s=\mu / 2$ for $\mu$ as in (*) is an element of $\pi^{\beta} \mathcal{N}_{k}^{e}(\overline{\mathbf{Q}})$, with $\beta$ and $e$ given there. (We write $e$ instead of $t$ employed there.) This result combined with (29.6) gives

$$
\begin{equation*}
\pi^{d_{0}-\gamma^{\prime}} \mathfrak{q}^{-1} \mathcal{Z}\left(\sigma_{0}, \mathbf{f}, \chi\right) f_{b}(z)=\sum_{a \in \mathcal{B}} \operatorname{vol}\left(\mathcal{D}_{a}\right) \sum_{i \in I_{a}}\left\langle h_{a, b, i}, f_{a}\right\rangle g_{a, b, i}(z) \tag{29.14}
\end{equation*}
$$

with $h_{a, b, i}, g_{a, b, i} \in \mathcal{N}_{h}^{e^{\prime}}(\overline{\mathbf{Q}})$, where $\gamma=\beta+\sum_{v \in \mathbf{a}} r_{v} q_{v}$ and $e^{\prime}$ is as in Lemma 29.2.
For the reason explained in $\S 28.2$ we have $\operatorname{vol}\left(\mathcal{D}_{a}\right)=\tau_{a} \pi^{d_{0}}$ with $0<\tau_{a} \in \overline{\mathbf{Q}}$. Let us now assume that $\varphi$ is anisotropic. Then $\mathcal{D}_{a}$ is compact, and so, by Proposition 15.7 (3), we can find $h_{a, b, i}^{\prime} \in \mathcal{S}_{h}^{\varphi}(\overline{\mathbf{Q}})$ such that $\tau_{a} h_{a, b, i}-h_{a, b, i}^{\prime}$ belongs to the set $\mathcal{T}_{\omega}^{p}$ in that proposition. Then the right-hand side of (29.14) is $\pi^{d_{0}} \sum_{a, i}\left\langle h_{a, b, i}^{\prime}, f_{a}\right\rangle g_{a, b, i}(z)$.

Let $w$ be a CM-point of $\mathfrak{Z}^{\varphi}$; let $\mathfrak{P}_{h}(w)$ denote a fixed element of $\mathbf{C}^{\times}$that represents the coset $\mathfrak{P}_{h}(w)$ defined by (11.17a). Then $g_{a, b, i}(w) \in \mathfrak{P}_{h}(w) \overline{\mathbf{Q}}$ by our definition of $\mathcal{N}_{h}^{e^{\prime}}(\overline{\mathbf{Q}})$. Therefore the last sum $\sum_{a, i}$ at $z=w$ can be written $\mathfrak{P}_{h}(w) \sum_{a}\left\langle h_{a, i}^{\prime \prime}, f_{a}\right\rangle$ with some $h_{a, i}^{\prime \prime} \in \mathcal{S}_{h}^{\varphi}(\overline{\mathbf{Q}})$. By means of the same technique as in Lemma 28.4 we can find an element $\mathbf{j}$ of $\mathcal{S}_{h}^{\varphi}\left(D^{\varphi}, \overline{\mathbf{Q}}\right)$ such that $\sum_{a}\left\langle h_{a, i}^{\prime \prime}, f_{a}\right\rangle=$ $\langle\mathbf{j}, \mathbf{f}\rangle$. Let $\mathbf{g}$ be the projection of $\mathbf{j}$ to $\mathcal{V}(\overline{\mathbf{Q}})$ with respect to the decomposition of $\mathcal{S}_{h}^{\varphi}\left(D^{\varphi}, \overline{\mathbf{Q}}\right)$ mentioned in §29.4. Observe that $\mathbf{g}$ depends on $b$ and $w$, but it is independent of $\mathbf{f}$. Writing $\mathbf{g}_{b, w}$ for $\mathbf{g}$, we thus have

$$
\begin{equation*}
\mathcal{Z}\left(\sigma_{0}, \mathbf{f}, \chi\right) \mathfrak{P}_{h}(w)^{-1} f_{b}(w)=\pi^{\gamma} \mathfrak{q}\left\langle\mathbf{g}_{b, w}, \mathbf{f}\right\rangle \tag{29.15}
\end{equation*}
$$

for every $\mathbf{f} \in \mathcal{V}(\overline{\mathbf{Q}})$. Now suppose $m_{0}>2 n$. Put $\varepsilon_{v}=0$ if $h_{v}+h_{v \rho}-m_{0} \in 2 \mathbf{Z}$ and $\varepsilon_{v}=1$ if $h_{v}+h_{v \rho}-m_{0} \notin 2 \mathbf{Z}$. Define $\ell \in \mathbf{Z}^{\mathbf{a}}$ as follows: $\ell_{v}=h_{v}-h_{v \rho}-\varepsilon_{v}$ if $r_{v}>0$ and $\ell_{v}=m_{0}-2 h_{v \rho}$ if $r_{v}=0$. By [S97, Lemma 11.14 (3)] we can find $\chi$ as in our theorem with this $\ell$. We can easily verify that (29.10a, b) are satified with $\sigma_{0}=m_{0} / 2$, and so we can consider (29.15) with this $\chi$ and $\mu=m_{0}$. By [S97, Proposition 20.4 (3)], $\mathcal{Z}\left(m_{0} / 2, \mathbf{f}, \chi\right) \neq 0$, since $m_{0}>2 n$. Therefore (29.15) shows that given $\mathbf{f}$, we can find $(b, w)$ so that $\left\langle\mathbf{g}_{b, w}, \mathbf{f}\right\rangle \neq 0$. Consequently the $\mathbf{g}_{b, w}$ for all $(b, w)$ span $\mathcal{V}$ over $\mathbf{C}$, and hence $\mathcal{V}(\overline{\mathbf{Q}})$ over $\overline{\mathbf{Q}}$. Now $\mathfrak{P}_{h}(w)^{-1} f_{b}(w) \in \overline{\mathbf{Q}}$, and hence from (29.15) we see that $\left\langle\mathbf{f}^{\prime}, \mathbf{f}\right\rangle \in \pi^{-\gamma} \mathfrak{q}^{-1} \mathcal{Z}\left(m_{0} / 2, \mathbf{f}, \chi\right) \overline{\mathbf{Q}}$ for every $\mathbf{f}, \mathbf{f}^{\prime} \in \mathcal{V}(\overline{\mathbf{Q}})$, from which (29.9) follows immediately.

Returning to an arbitrary $\sigma_{0}$ satisfying (29.10a, b, c), choose $(b, w)$ so that $f_{b}(w) \neq 0$. Dividing (29.15) by $\langle\mathbf{f}, \mathbf{f}\rangle$ and using the formula for $\beta$ in Theorem 17.12 (v), we obtain (29.12) when $\varphi$ is anisotropic.

Let us next assume that $\varphi$ is isotropic. Then we cannot use Proposition 15.7. Instead we use Lemma 29.3, which requires that $q_{v} \geq 0$ for every $v \in \mathbf{a}$. Now $\Lambda_{\mathrm{c}}^{n} H_{b, a}$ is $D_{r}$ of type (17.24) as we already noted, and its value at $s=\mu / 2$ is, by (17.27) or (17.30), of the form $\Delta_{\nu \mathbf{a}}^{p} R$ with $R \in \mathcal{M}_{\nu \mathbf{a}}$. Here $p=\left(m^{*}-\nu \mathbf{a}\right) / 2, \nu=\mu$ if $\mu \geq n$ and $\nu=2 n-\mu$ if $\mu<n ; R$ is of the form $D_{r}\left(\mathfrak{z}, \nu / 2 ; k^{\prime}, \chi, \mathfrak{c}\right)$ with $k^{\prime}$ such that $\left(k_{v}^{\prime}+k_{v \mathbf{a}}^{\prime}\right)_{v \in \mathbf{a}}=\nu \mathbf{a}$. Taking $\left(k^{\prime}, \nu\right)$ as $(k, \mu)$ of Theorem 17.12 (v), we find that $R \in \pi^{c} \mathcal{M}_{\nu \mathbf{a}}(\overline{\mathbf{Q}})$ with some $c \in \mathbf{Z}$, if we assume (29.10c) and (29.11b). By Lemma 29.3 we can replace $\Lambda_{\mathrm{c}}^{n}(\mu / 2, \chi)\left(A_{q}^{k} H_{b, a}\right)(z, w)$ by $\overline{T(z, w)}$ with $T$ of the form (29.7). Therefore we can repeat what we did in the anisotropic case, and obtain (29.14) under (29.11a), with $h_{a, b, i} \in \mathcal{M}_{h}^{\varphi}(\overline{\mathbf{Q}})$. Replacing these by elements of $\mathcal{S}_{h}^{\varphi}(\overline{\mathbf{Q}})$ by virtue of Theorem 27.14 (which is necessary only when $m \in \mathbf{Z a}$ ), we eventually obtain (29.15). Then we make a special choice $\ell_{v}=h_{v}-h_{v \rho}-\varepsilon_{v}$ with $\varepsilon_{v}=0$ or 1 . Since (29.11a) is satisfied by such an $\ell$, our proof of (29.9) is valid in the present case. Then (29.12) follows immediately from (29.15).
29.6. Theorem. Let the notation be as in Theorem 26.13 (see also [S97, Theorem 20.7]); let $q$ be as in §26.9. Let $\mathbf{f} \in \mathcal{V}(\overline{\mathbf{Q}})$ and $f \in \mathcal{S}_{k}^{\varphi}(\overline{\mathbf{Q}})$, where $\mathcal{V}(\overline{\mathbf{Q}})$
is defined within $\mathcal{S}_{k}^{\varphi}\left(D^{\varphi}\right)$ for a fixed system of eigenvalues $\{\lambda(\mathfrak{a})\}$ as in §29.3; put $m_{0}=\operatorname{Min}_{v \in \mathbf{a}^{\prime}}\left\{k_{v}+k_{v \rho}\right\}$; suppose that $m_{0}>2 n$ and $G_{\mathbf{a}}^{\varphi}$ is not compact.
(I) For $\nu$ as in Theorem 26.13 (i), the functions $\pi^{-\beta} E_{p}\left(z, \nu / 2 ; \mathbf{f}, \chi, D^{\psi}\right)$ is $\overline{\mathbf{Q}}$-rational, where $\beta=\sum_{v \in \mathbf{a}}\left(m_{v}-\nu\right) q / 2$, except when $\varphi$ is isotropic, $F=\mathbf{Q}$, $\nu=q+n+1$, and $\chi_{1}=\theta^{\nu}$; moreover, $\pi^{-\beta} E_{k}^{\psi, \varphi}(z, \nu / 2 ; f, \Gamma)$ with the same $\beta$ is $\overline{\mathbf{Q}}$-rational, except when $\varphi$ is isotropic, $F=\mathbf{Q}$, and $\nu=q+n+1$.
(II) For $\nu$ as in Theorem 26.13 (ii), the function $\pi^{-\alpha} \mathfrak{q}^{-1}\langle\mathbf{f}, \mathbf{f}\rangle^{-1} \mathcal{F}_{p}(z, \nu / 2 ; \mathbf{f}$, $\left.\chi, D^{\psi}\right)$ is $\overline{\mathbf{Q}}$-rational, where $\alpha=\sum_{v \in \mathbf{a}}\left(m_{v}+\nu-n-q+1\right)(n+q) / 2$ and $\mathfrak{q}=$ $p_{K}\left(\sum_{v \in \mathbf{a}} \ell_{v} \tau_{v}, \sum_{v \in \mathbf{a}} t_{v} \tau_{v}\right)$, provided the following two cases are excluded: (i) $0 \leq$ $\nu<q+n, \mathfrak{c}=\mathfrak{g}$, and $\chi_{1}=\theta^{\nu}$; (ii) $\varphi$ is isotropic, $F=\mathbf{Q}, \nu=q+n+1$, and $\chi_{1}=\theta^{\nu}$.

Proof. To prove (I), we evaluate (26.42) or its consequence (26.44) at $s=\nu / 2$. First suppose $\varphi$ is anisotropic. Then, by Proposition 15.7 (3), we can replace $h_{a i}$ of (26.44) by an element of $\mathcal{S}_{k}^{\varphi}(\overline{\mathbf{Q}})$. If $\varphi$ is isotropic, we express $H_{p, a}$ in the form $H_{p, a}=\Delta_{\nu \mathbf{a}}^{e} R$ with some $e$ and $R \in \mathcal{M}_{\nu \mathbf{a}}(\overline{\mathbf{Q}})$. This is possible by Theorem 17.12 (i) and (17.27). (Here we have to exclude Case (ii) of (II).) Then we apply Lemma 29.3 and (29.7) to $H_{p, a}$ and eventually find an expression of type (26.44) with $h_{a i} \in \mathcal{M}_{k}^{\varphi}(\overline{\mathbf{Q}})$. The whole procedure is similar to what was done in the proof of Theorem 29.5. Since $\operatorname{vol}\left(\mathcal{D}_{a}\right) \in \pi^{d_{0}} \overline{\mathbf{Q}}$, we eventually find, in both isotropic and anisotropic cases, that

$$
\begin{equation*}
\mathcal{Z}(\nu / 2, \mathbf{f}, \chi) E_{p}(z, \nu / 2)=\pi^{\alpha+\gamma_{\mathfrak{q}}} \sum_{a, i}\left\langle h_{a i}, f_{a}\right\rangle g_{a i}(z) \tag{29.16}
\end{equation*}
$$

with some $h_{a i} \in \mathcal{M}_{k}^{\varphi}(\overline{\mathbf{Q}})$ and $g_{a i} \in \mathcal{N}_{k}^{\psi, t}(\overline{\mathbf{Q}})$, where $\mathfrak{q}$ is as in Lemma 29.2, $\alpha=$ $(n+q) \sum_{v \in \mathbf{a}}\left(m_{v}-\nu\right) / 2$, and $\gamma=d n \nu-d n(n-1) / 2$. By the same technique as in the proof of Corollary 28.6 we can derive from (29.9) that $\left\langle h, f_{a}\right\rangle \in\langle\mathbf{f}, \mathbf{f}\rangle \overline{\mathbf{Q}}$ for every $h \in \mathcal{M}_{k}^{\varphi}(\overline{\mathbf{Q}})$ and every $a \in \mathcal{B}$. Also, $\mathcal{Z}(\nu / 2, \mathbf{f}, \chi) \neq 0$ for $\nu>2 n$ by [S97, Theorem 20.4 (3)]. Therefore, dividing (29.16) by $\langle\mathbf{f}, \mathbf{f}\rangle$ and employing (29.12) with $h=$ $k$, we obtain our assertion of (I) on $\pi^{-\beta} E_{p}\left(z, \ldots, D^{\varphi}\right)$, which combined with [S97, Proposition 20.10] proves the result about $\pi^{-\beta} E_{k}^{\psi, \varphi}(z, \ldots, \Gamma)$. Assertion (II) follows from (26.43) combined with Theorem 17.12 (v) in a similar way.
29.7. Theorem. The notation being as in $\S 29.4$, suppose that $\varphi$ is totally definite; let $\mathbf{f} \in \mathcal{V}(\overline{\mathbf{Q}})$.
(I) Let $\sigma_{0}$ be an element of $2^{-1} \mathbf{Z}$ such that

$$
\begin{equation*}
2 n-\left(2 h_{v \rho}+\ell_{v}\right) \leq 2 \sigma_{0} \leq 2 h_{v \rho}+\ell_{v} \text { and } 2 \sigma_{0}-\ell_{v} \in 2 \mathbf{Z} \text { for every } v \in \mathbf{a} \tag{29.17}
\end{equation*}
$$

Then, under (29.10c), $\mathcal{Z}\left(\sigma_{0}, \mathbf{f}, \chi\right) \in \pi^{\beta} p_{K}\left(\sum_{v \in \mathbf{a}}\left(2 h_{v \rho}+\ell_{v}\right) \tau_{v}, n \tau\right) \overline{\mathbf{Q}}$, where $\beta=$ $d n \sigma_{0}-d n(n-1) / 2+(n / 2) \sum_{v \in \mathbf{a}}\left(2 h_{v \rho}+\ell_{v}\right)$.
(II) Assertions (I) and (II) of Theorem 29.6 are true in the present case, if we replace $\langle\mathbf{f}, \mathbf{f}\rangle$ there by $p_{K}\left(\sum_{v \in \mathbf{a}} 2 k_{v \rho} \tau_{v}, n \tau\right)$ and put $t_{v}=n$ for every $v \in \mathbf{a}$.

Proof. Define $k \in \mathbf{Z}^{\mathbf{b}}$ by $k_{v \rho}=h_{v \rho}$ and $k_{v}=h_{v \rho}+\ell_{v}$ for $v \in \mathbf{a}$. Then $\ell=\left(k_{v}-k_{v \rho}\right)_{v \in \mathbf{a}}$. Moreover, by (29.2), $\mathcal{S}_{k}^{\varphi}\left(D^{\varphi}\right)=\mathcal{S}_{h}^{\varphi}\left(D^{\varphi}\right)$, and the Hecke operators on $\mathcal{S}_{h}^{\varphi}\left(D^{\varphi}\right)$ stay the same by the change of $h$ for $k$; also, $\mathcal{M}_{k}(\overline{\mathbf{Q}})=\mathcal{M}_{h}(\overline{\mathbf{Q}})$ by (26.41c). Therefore our calculations of [S97, §22.4 through §22.11] are valid with the present $k$. In particular, by [S97, (22.11.3)] we have

$$
\begin{equation*}
C^{\prime}(s) \mathcal{Z}(s, \mathbf{f}, \chi) f_{b}=\sum_{a \in \mathcal{B}} c_{a} f_{a} \Lambda_{\mathfrak{c}}^{n}(s, \chi)\left(H_{b, a}\right)^{\circ}(\mathbf{i}, \mathbf{i} ; s) \tag{*}
\end{equation*}
$$

where $c_{a} \in \overline{\mathbf{Q}}, C^{\prime}$ is as in [S97, (22.6.5)], and $\Lambda_{\mathrm{c}}^{n}$ is the function of (20.20) in Case UT. (See [S97, (20.3.3)]; also we may assume, without changing $\mathcal{Z}$, that the conductor of $\chi$ divides $\mathfrak{c}$ as noted in §26.10.) Put $\mu=2 \sigma_{0}$ and $M(\mathfrak{z}, s)=$ $\Lambda_{\mathfrak{c}}^{n}(s, \chi) H_{b, a}(\mathfrak{z}, s)$; notice that $k_{v}+k_{v \rho}=2 h_{v \rho}+\ell_{v}$. For the same reason as in the proof of Theorem 26.13, we can apply Theorem 17.12 (v) to $M(\mathfrak{z}, \mu / 2)$ to find that under (29.17) and (29.10c), $M(\mathfrak{z}, \mu / 2)$ belongs to $\pi^{\beta} \mathcal{N}_{k}^{p}(\overline{\mathbf{Q}})$ with $\beta=$ $(n / 2) \sum_{v \in \mathbf{a}}\left(m_{v}+\mu\right)-d n(n-1) / 2$ and some $p$. By $[\mathrm{S} 97,(22.6 .5)], C^{\prime}(\mu / 2) \in \overline{\mathbf{Q}}^{\times}$. By Lemma $26.12, M^{\circ}(\mathbf{i}, \mathbf{i} ; \mu / 2) \in \pi^{\beta} p_{K}\left(\sum_{v \in \mathbf{a}}\left(2 h_{v \rho}+\ell_{v}\right) \tau_{v}, n \tau\right) \overline{\mathbf{Q}}$. Take $b$ so that $f_{b} \neq 0$; then $f_{a} / f_{b} \in \overline{\mathbf{Q}}$ by (26.41c). Therefore we obtain (I) from ( ${ }^{*}$ ). Assertion (II) can be obtained from [S97, (22.11.2)] and Lemma 26.12 in a similar way.
29.8. Corollary. The notation being as in Theorems 29.6 and 29.7, suppose that $\mathbf{f}$ and $f$ are $\overline{\mathbf{Q}}$-rational, and $m=\nu \mathbf{a}$ with an integer $\nu \geq \operatorname{Max}\{2 n+$ $1, n+q\}$; suppose also that $\nu \neq n+q+1$ if $F=\mathbf{Q}$ and $\varphi$ is isotropic. Then $E_{k}^{\psi, \varphi}(z, \nu / 2 ; f, \Gamma)$ and $E_{p}\left(z, \nu / 2 ; \mathbf{f}, \chi, D^{\varphi}\right)$ are $\overline{\mathbf{Q}}$-rational.

This is an immediate consequence of those two theorems.
29.9. Proof of Theorem 27.16 (1) in Case UB. Our reasoning is similar to that of $\S 28.12$. We have $(W, \psi)=(V, \varphi) \oplus\left(H_{q}, \eta_{q}^{\prime}\right)$ and $(V, \varphi)=(Z, \zeta) \oplus\left(H_{r}, \eta_{r}^{\prime}\right)$ with $r=l-q$ as in $\S 26.4$ and Theorem 26.13; thus $\operatorname{dim}(W)=n+2 q$ with $n=\operatorname{dim}(V)=2 r+\operatorname{dim}(Z)$, and the lowest dimensional case is $\operatorname{dim}(W)=3$. We may assume that $\operatorname{dim}(Z)>0$, since if $Z=\{0\}$, then our group is reduced to Case UT. Also, there is no problem if $W=Z$, and so we assume that $l>0$. Since $\zeta$ is anisotropic, we have Theorems 29.5, 29.6, and 29.7 (without employing the map $\mathfrak{q}$ of Theorem 27.14) for the forms on $G^{\zeta}$ if $\mu>2 \operatorname{dim}(Z)$. Then we obtain Corollary 29.8 for $E_{k}^{\psi, \zeta}$, that is, Theorem 27.16 (1) for $r=0$. This establishes Theorem 27.16 in Case UB when $\operatorname{dim}(W)=3$. Now we make the induction assumption that the theorem is true for $G^{\varphi}$ with $\operatorname{dim}(V)<\operatorname{dim}(W)$. Let $0<r<l$ and let $\mu$ be as in Theorem 27.15. Our asumption guarantees the map $\mathfrak{q}$ for the forms on $G^{\varphi}$, since $\mu>2 n$. Then Theorems 29.5 and 29.6 are valid for the forms on $G^{\varphi}$ with such a $\mu$. Consequently we obtain Corollary 29.8 for $E_{k}^{\psi, \varphi}$, that is, Theorem 27.16 (1) for such $r$ and $\mu$. This completes the proof.
29.10. Remark. (A) In the above we stated the case of totally definite $\varphi$ separately, but we can state the results for anisotropic $\varphi$ uniformly, and view the totally definite case as a special case. The only point we must remember in that case is that $\mathcal{M}_{k}(\overline{\mathbf{Q}})=p_{K}\left(\sum_{v \in \mathbf{a}} k_{v \rho} \tau_{v}, n \tau\right) \overline{\mathbf{Q}}$ as stated in (26.41c), and consequently $\langle\mathbf{f}, \mathbf{f}\rangle \in p_{K}\left(\sum_{v \in \mathbf{a}} 2 k_{v \rho} \tau_{v}, n \tau\right) \overline{\mathbf{Q}}$. Also, the condition $m_{0}>2 n$ does not apply to the totally definite case.
(B) In Theorem 27.16 we stated a result of the type given in the above corollary, but the case $\nu=n+q \geq 2 n+1$ proved in the corollary is not included in that theorem. Also, the case in which $\nu=n+q+1 \geq 2 n+1, F=\mathbf{Q}$, and $\varphi$ is anisotropic is included in the above corollary, but not in Theorem 27.16.
(C) We assumed (29.11a) when $\varphi$ is isotropic, but probably this is unnecessary, since Lemma 29.3 is likely to be true without the condition that $q_{v} \geq 0$ for every $v$. At least the case $G^{\varphi}=U(\eta)$ is covered by Theorem 28.8. We can also handle some cases by means of Lemma 28.2 without assuming (29.11a). This method restricts $\sigma_{0}$ to a certain range, which is not large, but often nonempty. We leave the precise statement to the reader, as it is an easy exercise.

## APPENDIX

## A1. The series associated to a symmetric matrix and Gauss sums

A1.1. In this section we give some results necessary for the explicit calculations of the factors of automorphy of half-integral weight, and also prove the part of Theorem 16.2 concerning $A_{\zeta}^{1}$ and $f_{\zeta}^{1}$.

We consider $F_{v}$ with a fixed $v$ in $\mathbf{h}$ prime to 2 . For simplicity we drop the subscript $v$. Thus $F$ denotes a finite algebraic extension of $\mathbf{Q}_{p}$ with an odd prime number $p, \mathfrak{g}$ the integral closure of $\mathbf{Z}_{p}$ in $F$, and $\mathfrak{p}$ the maximal ideal of $\mathfrak{g}$. We denote by $\mathfrak{d}$ the different of $F$ relative to $\mathbf{Q}_{p}$ and by $\pi$ an unspecified prime element of $F$; also we put $q=[\mathfrak{g}: \mathfrak{p}]$. Let $\mathbf{e}_{v}$ be the $\mathbf{T}$-valued character of the additive group $F$ defined in $\S 1.6$. Here, to avoid a possible confusion, we use $\mathbf{e}_{v}$ without dropping the subscript $v$. Recall that $\mathfrak{d}^{-1}=\left\{a \in F \mid \mathbf{e}_{v}(a \mathfrak{g})=1\right\}$. We define the quadratic residue symbol $\left(\frac{c}{\mathfrak{p}}\right)$ for $c \in \mathfrak{g}$ as usual by the property that

$$
\begin{equation*}
1+\left(\frac{c}{\mathfrak{p}}\right)=\text { the number of } x(\bmod \mathfrak{p}) \text { such that } x^{2}-c \in \mathfrak{p} \tag{A1.1}
\end{equation*}
$$

We put then $\left(\frac{c}{\mathfrak{p}^{k}}\right)=\left(\frac{c}{\mathfrak{p}}\right)^{k}$ for every $k \in \mathbf{Z}$.
Given $a \in F^{\times}$such that $a \mathfrak{d}=\mathfrak{p}^{-m}$ with $0 \leq m \in \mathbf{Z}$, we put

$$
\begin{equation*}
\tau(a)=\sum_{x \in \mathfrak{g} / \mathfrak{p}^{m}} \mathbf{e}_{v}\left(a x^{2}\right) \tag{A1.2}
\end{equation*}
$$

This is well-defined.

## A1.2. Lemma.

$$
\tau(a)= \begin{cases}q^{m / 2} & \text { if } m \in 2 \mathbf{Z} \\ q^{(m-1) / 2} \sum_{y \in \mathfrak{g} / \mathfrak{p}}\left(\frac{y}{\mathfrak{p}}\right) \mathbf{e}_{v}\left(\pi^{m-1} a y\right) & \text { if } m \notin 2 \mathbf{Z}\end{cases}
$$

Proof. Assuming $m \geq 2$, we have

$$
\begin{aligned}
\tau(a) & =\sum_{y \in \mathfrak{g} / \mathfrak{p}^{m-1}} \sum_{z \in \mathfrak{g} / \mathfrak{p}} \mathbf{e}_{v}\left(a\left(y+\pi^{m-1} z\right)^{2}\right) \\
& =\sum_{y \in \mathfrak{g} / \mathfrak{p}^{m-1}} \mathbf{e}_{v}\left(a y^{2}\right) \sum_{z \in \mathfrak{g} / \mathfrak{p}} \mathbf{e}_{v}\left(2 a \pi^{m-1} y z\right)
\end{aligned}
$$

since $a \pi^{2 m-2} \in \mathfrak{d}^{-1}$. The last sum over $z$ is nonzero only if $y \in \mathfrak{p}$, in which case the sum is $q$. Thus $\tau(a)=q \tau\left(\pi^{2} a\right)$ if $m \geq 2$. Repeating this procedure, we
obtain $\tau(a)=q^{m / 2}$ if $m$ is even, since $\tau(a)=1$ if $m=0$. If $m$ is odd, we have $\tau(a)=q^{(m-1) / 2} \tau\left(\pi^{m-1} a\right)$. Putting $b=\pi^{m-1} a$ and employing (A1.1), we see that

$$
\tau(b)=\sum_{y \in \mathfrak{g} / \mathfrak{p}}\left\{1+\left(\frac{y}{\mathfrak{p}}\right)\right\} \mathbf{e}_{v}(b y)=\sum_{y \in \mathfrak{g} / \mathfrak{p}}\left(\frac{y}{\mathfrak{p}}\right) \mathbf{e}_{v}(b y)
$$

This completes the proof.
A1.3. Lemma. If $a \mathfrak{d}=\mathfrak{p}^{-m}$ as above, we have: (1) $\tau(a)^{2}=q^{m}\left(\frac{-1}{\mathfrak{p}}\right)^{m}$;
(2) $\tau(c a)=\tau(a)\left(\frac{c}{\mathfrak{p}}\right)^{m}$ if $c \in \mathfrak{g}^{\times}$.

Proof. The first equality for even $m$ and the second one for an arbitrary $m$ follow easily from Lemma A1.2. Put $\psi(y)=\left(\frac{y}{\mathfrak{p}}\right)$. Assuming $m=1$, we have

$$
\begin{aligned}
\tau(a) \overline{\tau(a)} & =\sum_{y \in(\mathfrak{g} / \mathfrak{p})^{\times}} \psi(y) \mathbf{e}_{v}(a y) \sum_{z \in(\mathfrak{g} / \mathfrak{p})^{\times}} \psi(z)^{-1} \mathbf{e}_{v}(-a z) \\
& =\sum_{y, z \in(\mathfrak{g} / \mathfrak{p})^{\times}} \psi\left(y z^{-1}\right) \mathbf{e}_{v}(a(y-z))=\sum_{x \in(\mathfrak{g} / \mathfrak{p})^{\times}} \psi(x) \sum_{z \in(\mathfrak{g} / \mathfrak{p})^{\times}} \mathbf{e}_{v}(a z(x-1)) .
\end{aligned}
$$

The last sum over $z$ is $q-1$ or -1 according as $x \equiv 1$ or $x \not \equiv 1(\bmod \mathfrak{p})$. Thus $|\tau(a)|^{2}=q-\sum_{x \in(\mathfrak{g} / \mathfrak{p}) \times} \psi(x)=q$. Clearly $\overline{\tau(a)}=\psi(-1) \tau(a)$, and hence $\tau(a)^{2}=\psi(-1) q$ if $m=1$. This together with Lemma A1.2 proves (1) for odd $m$.

A1.4. We now define symbols $S, L, \Lambda, \gamma(s)$, and $\omega(s)$ as follows:

$$
\begin{array}{ll}
S=S^{n}=\left\{\left.x \in F_{n}^{n}\right|^{t} x=x\right\}, & \\
L=\mathfrak{g}_{n}^{1}, \quad \Lambda=\Lambda^{n}=S^{n} \cap \mathfrak{g}_{n}^{n}, & \\
\gamma(s)=\int_{L} \mathbf{e}_{v}\left(x s \cdot{ }^{t} x / 2\right) d x & (s \in S) \\
\omega(s)=\nu(\delta s)^{1 / 2} \gamma(s) & (s \in S) \tag{A1.6}
\end{array}
$$

Here $d x$ is the Haar measure of $F_{n}^{1}$ such that $\int_{L} d x=1$, and $\delta$ is an element of $\mathfrak{g}$ such that $\mathfrak{d}=\delta \mathfrak{g} ; \nu()$ is defined by (1.16). (In the next section we use a different measure.) In Lemma A1.6 (4) below we shall show that $\omega(s)=\gamma(s) /|\gamma(s)|$, which is consistent with (16.6). If $n=1, s \in F^{\times}$, and $s \mathfrak{d}=\mathfrak{p}^{-m}$ with $0 \leq m \in \mathbf{Z}$, then

$$
\begin{equation*}
\tau(s / 2)=q^{m} \int_{\mathfrak{g}} \mathbf{e}_{v}\left(s x^{2} / 2\right) d x=q^{m} \gamma(s) \tag{A1.7}
\end{equation*}
$$

A1.5. Lemma. (1) Given $\sigma \in S^{n}$, there exists an element $\alpha$ of $G L_{n}(\mathfrak{g})$ such that $\alpha \sigma{ }^{t} \alpha$ is diagonal.
(2) If $\sigma \in S^{n} \cap G L_{n}(\mathfrak{g})$ and $n>1$, then $x \sigma \cdot{ }^{t} x=1$ for some $x \in \mathfrak{g}_{n}^{1}$.
(3) If $\sigma, \tau \in S^{n} \cap G L_{n}(\mathfrak{g})$ and $\operatorname{det}(\sigma) / \operatorname{det}(\tau)$ is a square in $\mathfrak{g} / \mathfrak{p}$, then $\sigma={ }^{t} \alpha \tau \alpha$ with some $\alpha \in G L_{n}(\mathfrak{g})$.

Proof. Since $2 x \sigma \cdot{ }^{t} y=(x+y) \sigma \cdot{ }^{t}(x+y)-x \sigma \cdot{ }^{t} x-y \sigma \cdot{ }^{t} y$ and $2 \notin \mathfrak{p}$, the ideal generated by $x \sigma \cdot{ }^{t} y$ for all $x, y \in L$ coincides with the ideal generated by $z \sigma \cdot{ }^{t} z$ for all $z \in L$. Thus, to prove (1), excluding the trivial case $\sigma=0$ and multiplying $\sigma$ by an element of $F^{\times}$, we may assume that this ideal is $\mathfrak{g}$. Take $z \in L$ so that $z \sigma \cdot{ }^{t} z \in \mathfrak{g}^{\times}$and put $M=\left\{y \in L \mid y \sigma \cdot{ }^{t} z=0\right\}$. Given $x \in L$, put $y=x-\left(x \sigma \cdot{ }^{t} z\right)\left(z \sigma \cdot{ }^{t} z\right)^{-1} z$. Then $y \in M$ and thus $L=\mathfrak{g} z \oplus M$. This means that
for some $\alpha \in G L_{n}(\mathfrak{g})$ we have $\alpha \sigma \cdot{ }^{t} \alpha=\operatorname{diag}\left[z \sigma \cdot{ }^{t} z, \tau\right]$ with $\tau \in \Lambda^{n-1}$. Applying induction to $\tau$, we obtain (1).

Clearly it is sufficient to prove (2) for diagonal $\sigma$ of size 2 . Thus put $\sigma=$ $\operatorname{diag}[a, d]$ with $a, d \in \mathfrak{g}^{\times}$. Suppose $d / a=-c^{2}$ with $c \in \mathfrak{g}^{\times}$. Then for $r, s \in \mathfrak{g}$ we have $a r^{2}+d s^{2}=a(r-c s)(r+c s)$. Since $2 \notin \mathfrak{p}$ we can find $r, s \in \mathfrak{g}$ so that $r+c s=1$ and $r-c s=a^{-1}$. Then $a r^{2}+d s^{2}=1$ as desired. If $-d / a$ is not a square, let $K=F(\xi)$ with $\xi^{2}=-d / a$. Then $K$ is an unramified quadratic extension of $F$, and the integral closure of $\mathfrak{g}$ in $K$ is $\mathfrak{g}[\xi]$; moreover $a \cdot N_{K / F}(r+s \xi)=a r^{2}+d s^{2}$. It is well- known that $N_{K / F}\left(\mathfrak{g}[\xi]^{\times}\right)=\mathfrak{g}^{\times}$, and hence $N_{K / F}(r+s \xi)=a^{-1}$ for some $r, s \in \mathfrak{g}$. This completes the proof of (2).

To prove (3), we first note that an element of $\mathfrak{g}^{\times}$is a square if and only if it is a square modulo $\mathfrak{p}$, since $2 \notin \mathfrak{p}$. Let $\sigma$ and $\tau$ be as in (3). Our assertion is trivial if $n=1$; thus we assume $n>1$. By (2) we have $x \sigma \cdot{ }^{t} x=1$ for some $x \in L$. Then the above proof of (1) shows that $\alpha \sigma \cdot{ }^{t} \alpha=\operatorname{diag}\left[1, \sigma^{\prime}\right]$ with some $\alpha \in G L_{n}(\mathfrak{g})$ and some $\sigma^{\prime} \in S^{n-1}$. Clearly $\sigma^{\prime} \in G L_{n-1}(\mathfrak{g})$. For the same reason we may assume that $\tau=\operatorname{diag}\left[1, \tau^{\prime}\right]$ with $\tau^{\prime} \in S^{n-1} \cap G L_{n-1}(\mathfrak{g})$. Then $\operatorname{det}\left(\sigma^{\prime}\right) / \operatorname{det}\left(\tau^{\prime}\right)$ is a square. Applying our induction to $\sigma^{\prime}$ and $\tau^{\prime}$, we obtain (3).

A1.6. Lemma. (1) $\gamma(c s)=\gamma(s)\left(\frac{c}{\nu_{0}(\delta s)}\right)$ if $c \in \mathfrak{g}^{\times}$, where $\nu_{0}$ is defined in §1.7;
(2) $\omega(s)^{2}=\left(\frac{-1}{\nu_{0}(\delta s)}\right) ;$ (3) $|\gamma(s)|=\nu(\delta s)^{-1 / 2}$;
(4) $\omega(s)=\gamma(s) /|\gamma(s)|$;
(5) $\gamma(s)=\gamma(s+b)$ and $\omega(s)=\omega(s+b)$ if $\delta b \in \Lambda$;
(6) $\gamma(-s)=\overline{\gamma(s)}$ and $\omega(-s)=\omega(s)^{-1}$.

Proof. By Lemma A1.5 (1) we may assume that $s=\operatorname{diag}\left[s_{1}, \ldots, s_{n}\right]$ with $s_{i} \in F$; we may also assume that $s_{i} \in \mathfrak{d}^{-1}$ if and only if $i>r$ with some $r \leq n$. Put $s_{i} \mathfrak{d}=\mathfrak{p}^{-m_{i}}$ with $0<m_{i} \in \mathbf{Z}$ for $i \leq r$ and $\lambda=\sum_{i=1}^{r} m_{i}$. Then $\nu_{0}(\delta s)=\mathfrak{p}^{\lambda}$, and

$$
\gamma(s)=\prod_{i=1}^{r} \int_{\mathfrak{g}} \mathbf{e}_{v}\left(s_{i} x^{2} / 2\right) d x=\prod_{i=1}^{r} q^{-m_{i}} \tau\left(s_{i} / 2\right)=\nu(\delta s)^{-1} \prod_{i=1}^{r} \tau\left(s_{i} / 2\right)
$$

by (A1.7). (If $r=0$, then $\nu_{0}(\delta s)=\mathfrak{g}$ and $\gamma(s)=1$.) Therefore, from Lemma A1.3 (2) we obtain (1). Also, by Lemma A1.3 (1) we have $\gamma(s)^{2}=\nu(\delta s)^{-1}\left(\frac{-1}{\nu_{0}(\delta s)}\right)$, which proves (2) and (3). Then $|\omega(s)|=1$, which implies (4). Clearly $\gamma(s)=$ $\gamma(s+b)$ and $\nu(\delta s)=\nu(\delta(s+b))$ if $\delta b \in \Lambda$. Therefore we obtain (5); the formulas of (6) are obvious.

A1.7. Given $\zeta \in \Lambda^{n}$, we define two infinite series $\alpha^{0}(\zeta, s)$ and $\alpha^{1}(\zeta, s)$ by

$$
\begin{align*}
& \alpha_{\zeta}^{0}(s)=\alpha^{0}(\zeta, s)=\sum_{\sigma \in S / \Lambda} \mathbf{e}_{v}^{n}\left(-\delta^{-1} \zeta \sigma\right) \nu(\sigma)^{-s}  \tag{A1.8}\\
& \alpha_{\zeta}^{1}(s)=\alpha^{1}(\zeta, s)=\sum_{\sigma \in S / \Lambda} \omega\left(\delta^{-1} \sigma\right) \mathbf{e}_{v}^{n}\left(-\delta^{-1} \zeta \sigma\right) \nu(\sigma)^{-s} . \tag{A1.9}
\end{align*}
$$

The sums are formally meaningful in view of Lemma A1.6 (5). These are the same as $\alpha_{v}^{i}$ of (16.7a, b). The series $\alpha_{\zeta}^{0}$ was investigated in [S97, $\left.\S \S 13-15\right]$. (In this section we assume $v \nmid 2$, and hence $T^{n}$ of [S97, (13.1.4)] coincides with the present $\Lambda^{n}$.) The purpose of the remaining part of this section is to determine $\alpha_{\zeta}^{1}$ as a rational function of $q^{-s}$ as stated in Theorem 16.2.

For positive integers $m$ and $n$ we put

$$
\begin{equation*}
\Lambda_{m}(n)=\mathfrak{p}^{-m} \Lambda / \Lambda \tag{A1.10}
\end{equation*}
$$

Define formal power series $A_{\zeta}^{i}(t)$ in an indeterminate $t$ as in Theorem 16.2, and using the symbol $e(\sigma)$ introduced there, define also $A_{\zeta}^{1, m}(t)$ by

$$
\begin{equation*}
A_{\zeta}^{1, m}(t)=\sum_{\sigma \in \Lambda_{m}(n)} \omega\left(\delta^{-1} \sigma\right) \mathbf{e}_{v}^{n}\left(-\delta^{-1} \zeta \sigma\right) t^{e(\sigma)} \quad(0<m \in \mathbf{Z}) \tag{A1.11}
\end{equation*}
$$

Clearly $A_{\zeta}^{1, m}(t)$ is a polynomial in $t$ and $\lim _{m \rightarrow \infty} A_{\zeta}^{1, m}(t)=A_{\zeta}^{1}(t)$ if $A_{\zeta}^{1}$ is convergent at $t$. In fact, $\alpha_{\zeta}^{0}$ is convergent for sufficiently large $\operatorname{Re}(s)$ as observed in [S97, p.104, lines 1-3], and hence the same is true for $\alpha_{\zeta}^{1}(s)$, since $\left|\omega\left(\delta^{-1} \sigma\right)\right|=1$.

To determine $\alpha_{\zeta}^{1}$, we need the number $N_{m}(\psi, \varphi)$ defined for $\varphi \in \Lambda^{n}$ and $\psi \in \Lambda^{h}$ defined by

$$
\begin{equation*}
N_{m}(\psi, \varphi)=\#\left\{\rho_{m}(x) \mid x \in \mathfrak{g}_{n}^{h},{ }^{t} x \psi x-\varphi \prec \mathfrak{p}^{m}\right\} \tag{A1.12}
\end{equation*}
$$

where we write $X \prec \mathfrak{p}^{m}$ if a matrix $X$ has entries in $\mathfrak{p}^{m}$ (as we did in §1.8) and $\rho_{m}$ is the natural map of $\mathfrak{g}_{n}^{h}$ to $\mathfrak{g}_{n}^{h} /\left(\mathfrak{p}^{m}\right)_{n}^{h}$ with any $h$ and $n$.

A1.8. Lemma. $N_{m}\left(1_{h}, \operatorname{diag}[1, \sigma]\right)=N_{m}\left(1_{h}, 1\right) N_{m}\left(1_{h-1}, \sigma\right)$ for every $\sigma \in$ $\Lambda^{n}$.

Proof. Suppose ${ }^{t} x x-\operatorname{diag}[1, \sigma] \prec \mathfrak{p}^{m}$ with $x \in \mathfrak{g}_{n+1}^{h}$; put $x=\left[\begin{array}{ll}u & y\end{array}\right]$ with $u \in \mathfrak{g}_{1}^{h}$ and $y \in \mathfrak{g}_{n}^{h}$. Then

$$
\begin{equation*}
{ }^{t} u u-1 \in \mathfrak{p}^{m}, \quad{ }^{t} u y \prec \mathfrak{p}^{m}, \quad{ }^{t} y y-\sigma \prec \mathfrak{p}^{m} . \tag{*}
\end{equation*}
$$

The number of $\rho_{m}(u)$ for $u$ satisfying the first relation is $N_{m}\left(1_{h}, 1\right)$. For each fixed $u$ we are going to show that the number of $\rho_{m}(y)$ for all $y$ satisfying the last two relations of $\left(^{*}\right)$ is $N_{m}\left(1_{h-1}, \sigma\right)$. Then we obtain our lemma. To do this, fix $u$, and take $a \in G L_{h}(\mathfrak{g})$ so that ${ }^{t} u a=\left[\begin{array}{ll}b & 0_{h-1}^{1}\end{array}\right]$ with $b \in \mathfrak{g}$. Let $c$ be the upper left entry of $a^{-1} \cdot{ }^{t} a^{-1}$ and $d$ the lower right submatrix of ${ }^{t} a a$ of size $h-1$; let $w=a^{-1} y$ and let $z$ (resp. $z^{\prime}$ ) be the lower $h-1$ rows (resp. the top row) of $w$. Then $b^{2} c-1 \in \mathfrak{p}^{m}$, and $y$ satisfies the last two relations of $\left(^{*}\right)$ if and only if $z^{\prime} \prec \mathfrak{p}^{m}$ and ${ }^{t} z d z-\sigma \prec \mathfrak{p}^{m}$. Thus the number of $\rho_{m}(y)$ equals $N_{m}(d, \sigma)$. Now $\operatorname{det}(d)=c \operatorname{det}(a)^{2}$ by Lemma 1.3 (1). Since $b^{2} c-1 \in \mathfrak{p}^{m}$, this means that $\operatorname{det}(d)$ is a square of a unit, and hence $d={ }^{t} e e$ with some $e \in G L_{h-1}(\mathfrak{g})$ by Lemma A1.5 (3). Thus $N_{m}(d, \sigma)=N_{m}\left(1_{h-1}, \sigma\right)$, which completes the proof.

A1.9. Proof of the part of Theorem 16.2 concerning $A_{\zeta}^{1}$ and $f_{\zeta}^{1}$. Let $0<k \in \mathbf{Z}$ and $\sigma \in \mathfrak{p}^{-m} \Lambda$. Taking the $k$-th power of (A1.5) and (A1.6), we obtain

$$
\nu(\sigma)^{-k / 2} \omega\left(\delta^{-1} \sigma\right)^{k}=q^{-m n k} \sum_{x \in\left(\mathfrak{g} / \mathfrak{p}^{m}\right)_{n}^{k}} \mathbf{e}_{v}^{k}\left(\delta^{-1} x \sigma \cdot{ }^{t} x / 2\right)
$$

where $\mathbf{e}_{v}^{k}(y)=\mathbf{e}_{v}(\operatorname{tr}(y))$ for a matrix $y$ of size $k$. Therefore, for $\zeta \in \Lambda$ we have

$$
\begin{aligned}
\sum_{\sigma \in \Lambda_{m}(n)} & \mathbf{e}_{v}^{n}\left(-\delta^{-1} \zeta \sigma\right) \omega\left(\delta^{-1} \sigma\right)^{k} \nu(\sigma)^{-k / 2} \\
& =q^{-m n k} \sum_{\sigma \in \Lambda_{m}(n)} \sum_{x \in\left(\mathfrak{g} / \mathfrak{p}^{m}\right)_{n}^{k}} \mathbf{e}_{v}^{n}\left(-\delta^{-1} \zeta \sigma\right) \mathbf{e}_{v}^{k}\left(\delta^{-1} x \sigma \cdot{ }^{t} x / 2\right) \\
& =q^{-m n k} \sum_{x \in\left(\mathfrak{g} / \mathfrak{p}^{m}\right)_{n}^{k}} \sum_{\sigma \in \Lambda_{m}(n)} \mathbf{e}_{v}^{n}\left((2 \delta)^{-1}\left({ }^{t} x x-2 \zeta\right) \sigma\right) \\
& =q^{m n(n+1) / 2-m n k} N_{m}\left(1_{k}, 2 \zeta\right)
\end{aligned}
$$

Let $A_{\zeta}^{0, m}(t)$ be the right-hand side of (A1.11) without the factor $\omega\left(\delta^{-1} \sigma\right)$. Put $\theta=\left(\frac{-1}{\mathfrak{p}}\right)$. By Lemma A1.6 (2), $\theta^{h e(\sigma)}=\omega\left(\delta^{-1} \sigma\right)^{2 h}$, and hence

$$
\begin{aligned}
A_{\zeta}^{0, m}\left(\theta^{h} q^{-h}\right) & =\sum_{\sigma \in \Lambda_{m}(n)} \mathbf{e}_{v}^{n}\left(-\delta^{-1} \zeta \sigma\right) \theta^{h e(\sigma)} q^{-h e(\sigma)} \\
& =\sum_{\sigma \in \Lambda_{m}(n)} \mathbf{e}_{v}^{n}\left(-\delta^{-1} \zeta \sigma\right) \omega\left(\delta^{-1} \sigma\right)^{2 h} \nu(\sigma)^{-h}
\end{aligned}
$$

Combining this with the above result, we obtain, for sufficiently large $h$,

$$
\begin{aligned}
A_{\zeta}^{0}\left(\theta^{h} q^{-h}\right) & =\lim _{m \rightarrow \infty} A_{\zeta}^{0, m}\left(\theta^{h} q^{-h}\right) \\
& =\lim _{m \rightarrow \infty} q^{m n(n+1) / 2-2 m n h} N_{m}\left(1_{2 h}, 2 \zeta\right)
\end{aligned}
$$

Similarly, taking $k=2 h+1$, we obtain
so that

$$
A_{\zeta}^{1, m}\left(\theta^{h} q^{-h-1 / 2}\right)=\sum_{\sigma \in \Lambda_{m}(n)} \mathbf{e}_{v}^{n}\left(-\delta^{-1} \zeta \sigma\right) \omega\left(\delta^{-1} \sigma\right)^{2 h+1} \nu(\sigma)^{-h-1 / 2},
$$

$$
A_{\zeta}^{1}\left(\theta^{h} q^{-h-1 / 2}\right)=\lim _{m \rightarrow \infty} q^{m n(n+1) / 2-m n(2 h+1)} N_{m}\left(1_{2 h+1}, 2 \zeta\right)
$$

Putting $\varepsilon=1 / 2$ and $\tau=\operatorname{diag}[\varepsilon, \zeta]$, by Lemma A1.8 we have

$$
\begin{aligned}
A_{\zeta}^{1}\left(\theta^{h} q^{-h-1 / 2}\right) & =\lim _{m \rightarrow \infty} \frac{q^{m(n+1)(n+2) / 2-m(n+1)(2 h+2)} N_{m}\left(1_{2 h+2}, 2 \tau\right)}{q^{m-m(2 h+2)} N_{m}\left(1_{2 h+2}, 1\right)} \\
& =A_{\tau}^{0}\left(\theta^{h+1} q^{-h-1}\right) / A_{\varepsilon}^{0}\left(\theta^{h+1} q^{-h-1}\right)
\end{aligned}
$$

Since this holds for infinitely many $h$, we have $A_{\zeta}^{1}(t)=A_{\tau}^{0}\left(\theta q^{-1 / 2} t\right) / A_{\varepsilon}^{0}\left(\theta q^{-1 / 2} t\right)$. Let $\zeta=\operatorname{diag}[\xi, 0]$ with $\xi \in \widetilde{S}_{v}^{r} \cap G L_{r}(F)$. By [S97, Theorem 13.6], $A_{\varepsilon}^{0}(t)=1-t$ and $A_{\tau}^{0}=f_{\tau} g_{\tau}$ with a polynomial $g_{\tau}$ with coefficients in $\mathbf{Z}$ whose constant term is 1 and a rational function $f_{\tau}$ given as follows:

$$
\begin{aligned}
& f_{\tau}(t)=\frac{(1-t) \prod_{i=1}^{[(n+1) / 2]}\left(1-q^{2 i} t^{2}\right)}{\left(1-\lambda q^{(2 n+1-r) / 2} t\right) \prod_{i=1}^{[(n-r) / 2]}\left(1-q^{2 n+2-r-2 i} t^{2}\right)} \quad \text { if } r \text { is odd }, \\
& f_{\tau}(t)=\frac{(1-t) \prod_{i=1}^{[(n+1) / 2]}\left(1-q^{2 i} t^{2}\right)}{\prod_{i=1}^{[(n-r+1) / 2]}\left(1-q^{2 n+3-r-2 i} t^{2}\right)} \quad \text { if } r \text { is even, }
\end{aligned}
$$

where $\lambda=\lambda(\operatorname{diag}[\varepsilon, \xi])$ with the symbol $\lambda()$ defined in $\S 16.1$. If $r$ is odd, we see that $\lambda=\theta \lambda(\xi)$; if further $\xi \in G L_{r}(\mathfrak{g})$, then $\operatorname{diag}[\varepsilon, \xi] \in G L_{r+1}(\mathfrak{g})$, and so $g_{\tau}=1$ as noted in [S97, Theorem 13.6]. If $\zeta=0$, then $\tau=\operatorname{diag}[\varepsilon, 0]$, and so $g_{\tau}=1$ by the same theorem. Therefore substituting $\theta q^{-1 / 2} t$ for $t$ in the formulas for $f_{\tau}$, we obtain our assertions concerning $A_{\zeta}^{1}$ and $f_{\zeta}^{1}$ of Theorem 16.2.

## A2. Metaplectic groups and factors of automorphy

A2.1. Our basic field is a totally real algebraic number field $F$ and we employ the symbols $\mathbf{a}, \mathbf{h}, \mathbf{v}$, and $\mathfrak{g}$ introduced in $\S 1.4$; we denote by $\mathfrak{d}$ the different of $F$ relative to $\mathbf{Q}$. For $x=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right] \in R_{2 n}^{2 n}$ with $a, b, c$, and $d$ in $R_{n}^{n}$, we put $a=a_{x}, b=b_{x}, c=c_{x}$, and $d=d_{x}$ if there is no fear of confusion.

With a positive integer $n$ we now put

$$
\begin{gathered}
X=F_{n}^{1}, \quad L=\mathfrak{g}_{n}^{1}, \quad L^{*}=\mathfrak{d}^{-1} L, \quad \eta=\eta_{n}=\left[\begin{array}{cc}
0 & -1_{n} \\
1_{n} & 0
\end{array}\right], \\
G=S p(n, F), \quad P=\left\{\alpha \in G \mid c_{\alpha}=0\right\}, \\
\Omega_{\mathbf{A}}=\left\{x \in G_{\mathbf{A}} \mid \operatorname{det}\left(c_{x}\right) \in F_{\mathbf{A}}^{\times}\right\}, \quad \Omega_{v}=\left\{x \in G_{v} \mid \operatorname{det}\left(c_{x}\right) \neq 0\right\} .
\end{gathered}
$$

(The present $G$ is $G$ of (3.26) in Case SP.) We let $G$ act on $X \times X=F_{2 n}^{1}$ by right multiplication. Then we can define the metaplectic groups $M p\left(X_{\mathbf{A}}\right)$ and $M p\left(X_{v}\right)$ for each $v \in \mathbf{v}$ in the sense of [W64] with respect to the alternating form $(x, y) \mapsto x \eta_{n} \cdot{ }^{t} y$ on $F_{2 n}^{1} \times F_{2 n}^{1}$. Recall that these groups, written $M_{\mathbf{A}}$ and $M_{v}$ for simplicity, are groups of unitary transformations on $L^{2}\left(X_{\mathbf{A}}\right)$ and on $L^{2}\left(X_{v}\right)$; there are exact sequences

$$
\begin{align*}
& 1 \longrightarrow \mathbf{T} \longrightarrow M_{\mathbf{A}} \longrightarrow G_{\mathbf{A}} \longrightarrow 1  \tag{A2.1a}\\
& 1 \longrightarrow \mathbf{T} \longrightarrow M_{v} \longrightarrow G_{v} \longrightarrow 1 \tag{A2.1b}
\end{align*} \quad(v \in \mathbf{v}) .
$$

We denote by pr the projection maps of $M_{\mathbf{A}}$ and $M_{v}$ to $G_{\mathbf{A}}$ and $G_{v}$. There is a natural lift $r: G \rightarrow M_{\mathbf{A}}$ by which we can consider $G$ a subgroup of $M_{\mathbf{A}}$. There are also two types of lifts

$$
\begin{equation*}
r_{P}: P_{\mathbf{A}} \rightarrow M_{\mathbf{A}}, \quad r_{\Omega}: \Omega_{\mathbf{A}} \rightarrow M_{\mathbf{A}} \tag{A2.2}
\end{equation*}
$$

which satisfy the formulas

$$
\begin{align*}
& {\left[r_{P}(\alpha) f\right](x)=\left|\operatorname{det}\left(a_{\alpha}\right)\right|_{\mathbf{A}}^{1 / 2} \mathbf{e}_{\mathbf{A}}\left(x a_{\alpha} \cdot{ }^{t} b_{\alpha} \cdot{ }^{t} x / 2\right) f\left(x a_{\alpha}\right) \text { if } \alpha \in P_{\mathbf{A}},}  \tag{A2.3a}\\
& r_{\Omega}(\alpha \beta \gamma)=r_{P}(\alpha) r_{\Omega}(\beta) r_{P}(\gamma) \text { if } \alpha, \gamma \in P_{\mathbf{A}} \quad \text { and } \beta \in \Omega_{\mathbf{A}}  \tag{A2.3b}\\
& {\left[r_{\Omega}(\beta) f\right](x)=\left|\operatorname{det}\left(c_{\beta}\right)\right|_{\mathbf{A}}^{1 / 2} \int_{X_{\mathbf{A}}} f\left(x a_{\beta}+y c_{\beta}\right) \mathbf{e}_{\mathbf{A}}\left(q_{\beta}(x, y)\right) d y \text { if } \beta \in \Omega_{\mathbf{A}},}  \tag{A2.3c}\\
& \quad q_{\beta}(x, y)=(1 / 2) x a_{\beta} \cdot{ }^{t} b_{\beta} \cdot{ }^{t} x+(1 / 2) y c_{\beta} \cdot{ }^{t} d_{\beta} \cdot{ }^{t} y+x b_{\beta} \cdot{ }^{t} c_{\beta} \cdot{ }^{t} y
\end{align*}
$$

Moreover $r_{P}=r$ on $P$ and $r_{\Omega}=r$ on $G \cap \Omega_{\mathbf{A}}$. There are also similar lifts of $P_{v}$ and $\Omega_{v}$ into $M_{v}$ given by the same formulas with the subscript $\mathbf{A}$ replaced by $v$. We denote these lifts also by $r_{P}$ and $r_{\Omega}$, since the distinction will be clear from the context. Here the measure on $X_{\mathbf{A}}$ is the $n$-fold product of the measure $\prod_{v \in \mathbf{a} \cup \mathbf{h}} d_{v} x$ on $F_{\mathbf{A}}, \int_{\mathfrak{g}_{v}} d_{v} x=N\left(\mathfrak{d}_{v}\right)^{-1 / 2}$ for $v \in \mathbf{h}, \int_{0}^{1} d_{v} x=1$ for $v \in \mathbf{a}$. We can let $S p(n, \mathbf{R})$ act on $\mathfrak{H}_{n}$ and define factors of automorphy $\mu(\alpha, z)$ for $\alpha \in S p(n, \mathbf{R})$ as in §3.3. We define a space $\mathcal{H}$ and a vector space $\mathcal{U}$ by

$$
\begin{equation*}
\mathcal{H}=\mathfrak{H}_{n}^{\mathbf{a}}, \quad \mathcal{U}=\left(\mathbf{C}^{n}\right)^{\mathbf{a}} \tag{A2.4}
\end{equation*}
$$

For $\sigma \in M_{\mathbf{A}}, \alpha=\operatorname{pr}(\sigma) \in G_{\mathbf{A}}$, and $z \in \mathcal{H}$ we put

$$
\begin{align*}
& \mu(\sigma, z)=\mu(\alpha, z)=\left(\mu\left(\alpha_{v}, z_{v}\right)\right)_{v \in \mathbf{a}}  \tag{A2.5a}\\
& j_{\sigma}(z)=j_{\alpha}(z)=\left(j\left(\alpha_{v}, z_{v}\right)\right)_{v \in \mathbf{a}}, \quad j\left(\alpha_{v}, z_{v}\right)=\operatorname{det}\left(\mu\left(\alpha_{v}, z_{v}\right)\right)  \tag{A2.5b}\\
& j_{\sigma}(z)^{\mathbf{a}}=j_{\alpha}(z)^{\mathbf{a}}=\prod_{v \in \mathbf{a}} j\left(\alpha_{v}, z_{v}\right) \tag{A2.5c}
\end{align*}
$$

We let $G_{v}=S p(n, \mathbf{R})$ act on $\mathbf{C}^{n} \times \mathfrak{H}_{n}$ by

$$
\begin{equation*}
\alpha(u, z)=\left({ }^{t} \mu(\alpha, z)^{-1} u, \alpha z\right) \tag{A2.6}
\end{equation*}
$$

for $\alpha \in G_{v}$, and let $G_{\mathbf{A}}$ act on $\mathcal{H}$ and $\mathcal{U} \times \mathcal{H}$ by

$$
\begin{equation*}
\alpha(z)=\left(\alpha_{v} z_{v}\right)_{v \in \mathbf{a}}, \quad \alpha(u, z)=\left({ }^{t} \mu\left(\alpha_{v}, z_{v}\right)^{-1} u_{v}, \alpha_{v} z_{v}\right)_{v \in \mathbf{a}} \tag{A2.7}
\end{equation*}
$$

for $z \in \mathcal{H}$ and $\alpha \in G_{\mathbf{A}}$, ignoring $\alpha_{\mathrm{h}}$. We define the action of an element $\sigma$ of $M_{\mathbf{A}}$ (resp. $M_{v}$ ) on $\mathcal{H}$ and $\mathcal{U} \times \mathcal{H}$ (resp. $\mathfrak{H}_{n}$ and $\mathbf{C}^{n} \times \mathfrak{H}_{n}$ ) to be the same as that of $\operatorname{pr}(\sigma)$.

We now define a $\mathbf{C}$-valued function $\varphi(x ; u, z)$ for $x \in X_{\mathbf{a}}, z \in \mathcal{H}$, and $u \in \mathcal{U}$ by

$$
\begin{gather*}
\varphi(x ; u, z)=\prod_{v \in \mathbf{a}} \varphi_{v}\left(x_{v} ; u_{v}, z_{v}\right)  \tag{A2.8a}\\
\varphi_{v}\left(x_{v} ; u_{v}, z_{v}\right)=\mathbf{e}\left((1 / 2)^{t} u(z-\bar{z})^{-1} u+(1 / 2) x z \cdot{ }^{t} x+x u\right), \tag{A2.8b}
\end{gather*}
$$

where the subscript $v$ is suppressed on the right-hand side of the last formula.
A2.2. Lemma. Let $M^{\prime}$ be the group formed by all the couples $(\alpha, g)$ with $\alpha \in$ $S p(n, \mathbf{R})$ and a holomorphic function $g$ on $\mathfrak{H}_{n}$ such that $g(z)^{2}=t \cdot j_{\alpha}(z)$ with $t \in \mathbf{T}$, the law of composition being $(\alpha, g)\left(\alpha^{\prime}, g^{\prime}\right)=\left(\alpha \alpha^{\prime}, g\left(\alpha^{\prime}(z)\right) g^{\prime}(z)\right)$. Then for each $v \in \mathbf{a}, M_{v}$ is isomorphic to $M^{\prime}$ via the map $\xi \mapsto\left(\operatorname{pr}(\xi), g_{\xi}\right) \in M^{\prime}$ for $\xi \in M_{v}$ with $g_{\xi}$ determined by

$$
\begin{equation*}
\left(\xi \varphi_{v}\right)(x ; u, z)=g_{\xi}(z)^{-1} \varphi_{v}(x ; \xi(u, z)) \quad\left(\xi \in M_{v}\right) \tag{A2.9}
\end{equation*}
$$

In particular $g_{\xi}(z)=\operatorname{det}(-i z)^{1 / 2}$ if $\xi=r_{\Omega}\left(\eta_{n}\right)$.
Proof. Let $\tau=r_{P}(\alpha)$ with $\alpha \in P$. Then from (A2.3a) we can easily derive that $\left(\tau \varphi_{v}\right)(x ; u, z)=g_{\tau}^{-1} \varphi_{v}(x ; \tau(u, z))$ with $g_{\tau}=\left|\operatorname{det}\left(d_{\alpha}\right)\right|_{v}^{1 / 2}$. Also. if $\sigma=$ $r_{\Omega}\left(\eta_{n}\right)$, formula (A2.3c) together with an easy calculation shows that

$$
\begin{align*}
\left(\sigma \varphi_{v}\right)(x ; u, z) & =\int_{\mathbf{R}_{n}^{1}} \varphi_{v}(y ; u, z) \mathbf{e}\left(-x \cdot{ }^{t} y\right) d y  \tag{A2.10}\\
& =\operatorname{det}(-i z)^{-1 / 2} \varphi_{v}(x ; \sigma(u, z))
\end{align*}
$$

Now from (A2.1b) and Lemma 7.5 we see that $M_{v}$ is generated by $r_{P}\left(P_{v}\right), r_{\Omega}\left(\eta_{n}\right)$, and $\mathbf{T}$. Therefore we have (A2.9) with a certain $g_{\xi}$ such that $g_{\xi}(z)^{2} / j_{\alpha}(z) \in \mathbf{T}$. Then $\xi \mapsto\left(\operatorname{pr}(\xi), g_{\xi}\right)$ defines a homomorphism of $M_{v}$ into $M^{\prime}$. In particular, an element $t$ of $\mathbf{T}$ is mapped to $\left(1, t^{-1}\right)$. Therefore we see that the map is injective. Given $(\alpha, g) \in M^{\prime}$, take $\xi \in M_{v}$ so that $\operatorname{pr}(\xi)=\alpha$. Then $g_{\xi}^{2} / g^{2} \in \mathbf{T}$, so that $g_{\xi}=t g$ with $t \in \mathrm{~T}$, and so $(\alpha, g)$ is the image of $t \xi$. Thus our map is surjective. The last assertion of our lemma follows from (A2.10).

Thus, writing $g(\xi, z)$ for $g_{\xi}(z)$, we have $g\left(\xi \xi^{\prime}, z\right)=g\left(\xi, \xi^{\prime} z\right) g\left(\xi^{\prime}, z\right)$ for $\xi, \xi^{\prime} \in$ $M_{v}$. Moreover, we have shown that

$$
\begin{equation*}
g\left(r_{P}(\alpha), z\right)=\left|\operatorname{det}\left(d_{\alpha}\right)\right|_{v}^{1 / 2} \quad \text { if } \quad \alpha \in P_{v}, v \in \mathbf{a} . \tag{A2.11}
\end{equation*}
$$

A2.3. For two fractional ideals $\mathfrak{x}$ and $\mathfrak{y}$ in $F$ such that $\mathfrak{x y} \subset \mathfrak{g}$ we define a subgroup $D[\mathfrak{x}, \mathfrak{y}]$ of $G_{\mathbf{A}}$ by (16.20a) and put $D_{v}[\mathfrak{x}, \mathfrak{y}]=G_{v} \cap D[\mathfrak{x}, \mathfrak{y}]$; we define also a subgroup $C^{\theta}$ of $G_{\mathbf{A}}$, which may be called the "theta-subgroup," by

$$
\begin{align*}
C^{\theta}= & G_{\mathbf{a}} \prod_{v \in \mathbf{h}} C_{v}^{\theta}  \tag{A2.12a}\\
C_{v}^{\theta}=\left\{\xi \in D_{v}\left[\mathfrak{d}^{-1}, \mathfrak{d}\right] \mid\right. & \chi_{v}((x, y) \xi)=\chi_{v}(x, y)  \tag{A2.12b}\\
& \text { for every } \left.x \in L_{v} \text { and } y \in L_{v}^{*}\right\} \quad \text { if } v \in \mathbf{h}
\end{align*}
$$

where $\chi_{v}(x, y)=\mathbf{e}_{v}\left(x \cdot{ }^{t} y / 2\right)$ for $x, y \in\left(F_{v}\right)_{n}^{1}$. Then it can be shown that

$$
\begin{align*}
C_{v}^{\theta}=\left\{\xi \in D_{v}\left[\mathfrak{d}^{-1}, \mathfrak{d}\right] \mid\right. & \left(a_{\xi} \cdot{ }^{t} b_{\xi}\right)_{i i} \in 2 \mathfrak{d}_{v}^{-1} \quad \text { and }  \tag{A2.13}\\
& \left.\left(c_{\xi} \cdot{ }^{t} d_{\xi}\right)_{i i} \in 2 \mathfrak{d}_{v} \quad \text { for } 1 \leqq i \leqq n\right\},
\end{align*}
$$

where $\alpha_{i i}$ denotes the ( $i, i$ )-entry of $\alpha$. We see from (A2.12b) that $C_{v}^{\theta}$ is indeed a group. Also we easily see that

$$
\begin{equation*}
C_{v}^{\theta} \supset D_{v}\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right] \cup D_{v}\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right] \varepsilon_{v} \tag{A2.14}
\end{equation*}
$$

with $\varepsilon \in G_{\mathbf{A}}$ given by

$$
\varepsilon_{\mathbf{a}}=1_{2 n}, \quad \varepsilon_{v}=\left[\begin{array}{cc}
0 & -\delta_{v}^{-1} 1_{n}  \tag{A2.15}\\
\delta_{v} 1_{n} & 0
\end{array}\right] \text { for } \quad v \in \mathbf{h}
$$

where $\delta$ is an arbitrarily fixed element of $F_{\mathbf{h}}^{\times}$such that $\mathfrak{d}=\delta \mathfrak{g}$. Since $C_{v}^{\theta}$ coincides with $\operatorname{Ps}\left(X_{v}, L_{v}\right)$ of [W64, $\left.\mathrm{n}^{\circ} 36\right]$, we obtain a lift

$$
\begin{equation*}
r_{v}: C_{v}^{\theta} \rightarrow M_{v} \tag{A2.16}
\end{equation*}
$$

which is written $\mathbf{r}_{L}^{\prime}$ there.
We are going to define a factor of automorphy, $h(\sigma, z)$, of weight $\mathbf{a} / 2$ for $z \in \mathcal{H}$ and $\sigma$ in the set

$$
\begin{equation*}
\mathfrak{M}=\left\{\sigma \in M_{\mathbf{A}} \mid \operatorname{pr}(\sigma) \in P_{\mathbf{A}} C^{\theta}\right\} \tag{A2.17}
\end{equation*}
$$

Clearly $r_{P}\left(P_{\mathbf{A}}\right) \subset \mathfrak{M}$. Notice that $\eta \in G \cap \mathfrak{M}$ since $\eta_{\mathbf{h}}=\operatorname{diag}\left[\delta 1_{n}, \delta^{-1} 1_{n}\right] \varepsilon \in P_{\mathbf{A}} C^{\theta}$.
We denote by $\mathcal{S}\left(X_{\mathrm{h}}\right)$ and $\mathcal{S}\left(X_{\mathbf{A}}\right)$ the Schwartz-Bruhat spaces of $X_{\mathrm{h}}$ and $X_{\mathrm{A}}$. We shall often view an element $\ell$ of $\mathcal{S}\left(X_{\mathrm{h}}\right)$ as a function on $X$ by restricting $\ell$ to the image of $X$ in $X_{\mathrm{h}}$ (see §1.6). Given $\ell \in \mathcal{S}\left(X_{\mathrm{h}}\right)$, we put

$$
\begin{equation*}
\ell_{\mathbf{A}}(x ; u, z)=\ell\left(x_{\mathbf{h}}\right) \varphi\left(x_{\mathbf{a}} ; u, z\right) \quad \text { for } x \in X_{\mathbf{A}}, z \in \mathcal{H}, u \in \mathcal{U} \tag{A2.18}
\end{equation*}
$$

For fixed $z$ and $u$ we view $\ell_{\mathbf{A}}$ as an element of $\mathcal{S}\left(X_{\mathbf{A}}\right)$, so that $\sigma \ell_{\mathbf{A}}$ for $\sigma \in M_{\mathbf{A}}$ is meaningful. Now there is another action of $\mathfrak{M}$ on $\mathcal{S}\left(X_{\mathbf{h}}\right)$ as follows:

A2.4. Theorem. For every $\sigma \in \mathfrak{M}$ we can define its action on $\mathcal{S}\left(X_{\mathbf{h}}\right)$ which is $a \mathbf{C}$-linear automorphism, written $\ell \mapsto^{\sigma} \ell$ for $\ell \in \mathcal{S}\left(X_{\mathbf{h}}\right)$, and also a holomorphic function $h(\sigma, z)$ of $z \in \mathcal{H}$ by the formula

$$
\begin{equation*}
\left(\sigma \ell_{\mathbf{A}}\right)(x ; u, z)=h(\sigma, z)^{-1}\left({ }^{\sigma} \ell\right)_{\mathbf{A}}(x ; \sigma(u, z)) . \tag{1}
\end{equation*}
$$

Moreover, this action and $h$ have the following properties:
(2) $h(\sigma, z)^{2}=\zeta j_{\sigma}(z)^{\mathbf{a}}$ with $\zeta \in \mathbf{T}$.
(3) $h\left(t \cdot r_{P}(\gamma), z\right)=t^{-1}\left|\operatorname{det}\left(d_{\gamma}\right)_{\mathbf{a}}\right|_{\mathbf{A}}^{\mathbf{T}} \quad$ if $t \in \mathbf{T}$ and $\gamma \in P_{\mathbf{A}}$.
(4) $h(\rho \sigma \tau, z)=h(\rho, z) h(\sigma, \tau z) h(\tau, z)$ and $(\rho \sigma \tau) \ell={ }^{\rho}\left({ }^{\sigma}\left({ }^{\tau} \ell\right)\right)$ if $\operatorname{pr}(\rho) \in P_{\mathbf{A}}$ and $\operatorname{pr}(\tau) \in C^{\theta}$; in particular, $h(-\sigma, z)=h(\sigma, z)$.
(5) ${ }^{\sigma} \ell$ depends only on $\ell$ and $\operatorname{pr}(\sigma)_{\mathbf{h}}$.
(6) $\left\{\left.\sigma \in \mathfrak{M}\right|^{\sigma} \ell=\ell\right\}$ contains an open subgroup of $M_{\mathbf{A}}$ for every $\ell \in \mathcal{S}\left(X_{\mathrm{h}}\right)$. In particular, if $\ell$ is the characteristic function of $\prod_{v \in \mathbf{h}} L_{v}$, then ${ }^{\sigma} \ell=\ell$ for every $\sigma$ such that $\operatorname{pr}(\sigma) \in C^{\theta}$.
(7) If $\ell(x)=\prod_{v \in \mathbf{h}} \ell_{v}\left(x_{v}\right)$ with $\ell_{v} \in \mathcal{S}\left(X_{v}\right), \ell_{v}$ is the characteristic function of $L_{v}$ for almost all $v$, and $\sigma=r_{\Omega}(\tau), \tau \in P_{\mathbf{A}} C^{*}$ with $C^{*}=\left\{\alpha \in C^{\theta} \mid L_{v}^{*}\left(c_{\alpha}\right)_{v}=\right.$ $L_{v}$ for every $\left.v \in \mathbf{h}\right\}$, then ${ }^{\sigma} \ell(x)=\prod_{v \in \mathbf{h}}\left[r_{\Omega}\left(\tau_{v}\right) \ell_{v}\right]\left(x_{v}\right)$.
(8) If $\ell$ is as in (7) and $\operatorname{pr}(\sigma)=\beta \alpha$ with $\beta \in P_{\mathbf{A}}$ and $\alpha \in C^{\theta}$, then ${ }^{\sigma} \ell(x)=$ $\prod_{v \in \mathbf{h}}\left[r_{P}\left(\beta_{v}\right) r_{v}\left(\alpha_{v}\right) \ell_{v}\right]\left(x_{v}\right)$, where $r_{v}$ is the lift of (A2.16).

Proof. Let $\tau \in M_{\mathrm{A}}$ and $\alpha=\operatorname{pr}(\tau)$; suppose $\alpha \in C^{\theta}$. Take $\xi_{v} \in M_{v}$ for each $v \in$ a so that $\operatorname{pr}\left(\xi_{v}\right)=\alpha_{v}$. Then we can define an element $\gamma$ of $M_{\mathbf{A}}$ such that $\operatorname{pr}(\gamma)=\alpha$ and

$$
\begin{equation*}
\left(\gamma \ell_{\mathbf{A}}\right)(x ; u, z)=\prod_{v \in \mathbf{h}}\left[r_{v}\left(\alpha_{v}\right) \ell_{v}\right]\left(x_{v}\right) \prod_{v \in \mathbf{a}}\left(\xi_{v} \varphi_{v}\right)\left(x_{v} ; u_{v}, z_{v}\right) \tag{A2.19}
\end{equation*}
$$

for $\ell=\prod_{v \in \mathrm{~h}} \ell_{v}$ as in (7) (see [W64, $\left.\mathrm{n}^{\circ} 38\right]$ ). Then $\tau=\zeta \gamma$ with $\zeta \in \mathbf{T}$. Now every element $\sigma$ of $\mathfrak{M}$ can be written $\sigma=r_{P}(\beta) \tau$ with such a $\tau$ and $\beta \in P_{\mathbf{A}}$. Applying $r_{P}(\beta)$ to (A2.19), we obtain

$$
\begin{equation*}
\left(\sigma \ell_{\mathbf{A}}\right)(x ; u, z)=\zeta \ell^{\prime}\left(x_{\mathbf{h}}\right) \prod_{v \in \mathbf{a}}\left[r_{P}\left(\beta_{v}\right) \xi_{v} \varphi_{v}\right]\left(x_{v} ; u_{v}, z_{v}\right) \tag{A2.20}
\end{equation*}
$$

where $\ell^{\prime}=\prod_{v \in \mathbf{h}} r_{P}\left(\beta_{v}\right) r_{v}\left(\alpha_{v}\right) \ell_{v}$. By (A2.9) we can write (A2.20) in the form

$$
\begin{equation*}
\left(\sigma \ell_{\mathbf{A}}\right)(x ; u, z)=h(z)^{-1} \ell_{\mathbf{A}}^{\prime}(x ; \sigma(u, z)) \tag{A2.21}
\end{equation*}
$$

with $h(z)=\zeta^{-1} \prod_{v \in \mathbf{a}} g\left(r_{P}\left(\beta_{v}\right) \xi_{v}, z_{v}\right)$. We have assumed that $\ell=\prod_{v \in \mathbf{h}} \ell_{v}$, but clearly $\ell \mapsto \ell^{\prime}$ can be extended to a C-linear automorphism of $\mathcal{S}\left(X_{\mathbf{h}}\right)$, which we write again $\ell \mapsto \ell^{\prime}$. Then (A2.21) holds for every $\ell \in \mathcal{S}\left(X_{\mathbf{h}}\right)$ with the same $h(z)$. We put $h(\sigma, z)=h(z)$ and ${ }^{\sigma} \ell=\ell^{\prime}$. To show that these are independent of the choice of $\beta$ and $\tau$, take $\ell$ to be the characteristic function $\lambda$ of $\prod_{v \in \mathbf{h}} L_{v}$, and recall that $r_{v}\left(C_{v}^{\theta}\right) \lambda_{v}=\lambda_{v}$ (see [W64, $\left.\mathrm{n}^{\circ} 21\right]$ ). Now (A2.3a) shows that $\left(r_{P}\left(\beta_{v}\right) \lambda_{v}\right)(0)=$ $\left|\operatorname{det}\left(d_{\beta}\right)\right|_{v}^{-1 / 2}$. Therefore, putting $x=0$ and $u=0$ in (A2.21), we obtain

$$
\begin{equation*}
\left(\sigma \lambda_{\mathbf{A}}\right)(0 ; 0, z)=\left|\operatorname{det}\left(d_{\beta}\right)_{\mathbf{h}}\right|_{\mathbf{A}}^{-1 / 2} \cdot h(z)^{-1} \tag{A2.22}
\end{equation*}
$$

Since $\operatorname{pr}(\sigma) \in \beta D\left[\mathfrak{d}^{-1}, \mathfrak{d}\right]$, from (1.19) we obtain $\left|\operatorname{det}\left(d_{\beta}\right)_{\mathbf{h}}\right|_{\mathbf{A}}^{-1}=N\left(\mathrm{il}_{\mathfrak{d}}(\operatorname{pr}(\sigma))\right)$, which depends only on $\operatorname{pr}(\sigma)$. Thus $h(z)$ is determined by $\sigma$, and consequently formula (1) is established with ${ }^{\sigma} \ell$ well-defined. Clearly (2), (4), (5), and (8) follow easily from our definition of $h(\sigma, z)$ and ${ }^{\sigma} \ell$; (3) follows from (A2.11) if we take $\xi_{v}=1$ and $\alpha=1$; the first part of (6) can be derived from the fact that $\left\{\gamma \in C_{v}^{\theta} \mid r_{v}(\gamma) \ell_{v}=\ell_{v}\right\}$ is open for every $\ell_{v} \in \mathcal{S}\left(X_{v}\right)$ (see [W64, $\mathrm{n}^{\circ} 21$, $\left.\mathrm{n}^{\circ} 36\right]$ ). Finally let $\sigma=r_{\Omega}(\tau)$ with $\tau=\beta \alpha$ with $\beta \in P_{\mathbf{A}}$ and $\alpha \in C^{*}$. Then $r_{\Omega}\left(\tau_{v}\right)=r_{P}\left(\beta_{v}\right) r_{\Omega}\left(\alpha_{v}\right) ;$ moreover $r_{\Omega}\left(\alpha_{v}\right)=r_{v}\left(\alpha_{v}\right)$ by [W64, p.168, last line]. Therefore (7) follows from (8). The second half of (6) follows also from (8).

Given $\ell \in \mathcal{S}\left(X_{\mathbf{h}}\right)$, we define a theta function $\theta(u, z ; \ell)$ for $(u, z) \in \mathcal{U} \times \mathcal{H}$ by

$$
\begin{equation*}
\theta(u, z ; \ell)=\sum_{\xi \in X} \ell_{\mathbf{A}}(\xi ; u, z) \tag{A2.23}
\end{equation*}
$$

A2.5. Proposition. For every $\alpha \in G \cap \mathfrak{M}$ we have

$$
\theta\left(\alpha(u, z) ;{ }^{\alpha} \ell\right)=h(\alpha, z) \theta(u, z ; \ell) .
$$

Moreover, for every $\beta \in G$ and $\xi \in \prod_{v \in \mathbf{a}} M_{v}$ such that $\operatorname{pr}(\xi)=\beta_{\mathbf{a}}$, there is a $\mathbf{C}$-linear automorphism $\ell \mapsto \ell^{\prime}$ of $\mathcal{S}\left(X_{\mathbf{h}}\right)$ such that

$$
\theta\left(\beta(u, z) ; \ell^{\prime}\right)=\theta(u, z ; \ell) \prod_{v \in \mathbf{a}} g\left(\xi_{v}, z_{v}\right)
$$

Proof. By virtue of [W64, Theorem 4 or 6] we have $\sum_{\xi \in X}\left(\alpha \ell_{\mathbf{A}}\right)(\xi ; u, z)=$ $\sum_{\xi \in X} \ell_{\mathbf{A}}(\xi ; u, z)$ for every $\alpha \in G$, which combined with (1) of Theorem A2.4 proves the first assertion. Now, by Lemma $7.5, G$ is generated by $G \cap \mathfrak{M}$, since $G \cap \mathfrak{M}$ contains $\eta$ and $P$. Therefore the second assertion follows from the first one.

To state our formulas on $h(\sigma, z)$, we need the symbols $\gamma(s)$ and $\omega(s)$ of (16.5) and (16.6), as well as $\mathrm{il}_{3}$ of (1.19).

A2.6. Lemma. Let $\sigma \in \mathfrak{M}$; if $\sigma$ or $\sigma \eta^{-1}$ belongs to $r_{\Omega}\left(\Omega_{\mathbf{A}} \cap P_{\mathbf{A}} C^{\theta}\right)$, then $h(\sigma, z)$ is completely determined by Theorem A2.4 (2) and the following formulas:

$$
\begin{aligned}
& \lim _{r \rightarrow \infty} h(\sigma, r \mathbf{i}) /|h(\sigma, r \mathbf{i})|=\omega\left(-c_{\alpha}^{-1} d_{\alpha}\right) \text { if } \sigma=r_{\Omega}(\alpha) \text { with } \alpha \in \Omega_{\mathbf{A}} \cap P_{\mathbf{A}} C^{\theta} \\
& \lim _{r \rightarrow 0} h(\sigma, r \mathbf{i}) /|h(\sigma, r \mathbf{i})|=\omega\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right) \text { if } \sigma=r_{\Omega}\left(\alpha \eta^{-1}\right) \eta \text { with } \alpha \in \Omega_{\mathbf{A}} \eta \cap P_{\mathbf{A}} C^{\theta}
\end{aligned}
$$

where $\mathbf{i}$ is the origin of $\mathcal{H}$ defined by (16.21). In particular

$$
\begin{equation*}
h(\eta, z)=\prod_{v \in \mathbf{a}} \operatorname{det}\left(-i z_{v}\right)^{1 / 2} \tag{A2.24}
\end{equation*}
$$

where $\operatorname{det}\left(-i z_{v}\right)^{1 / 2}$ is chosen so that it is positive when $\operatorname{Re}\left(z_{v}\right)=0$.
In order to speak of $\omega(s)$, we need to know $\gamma(s) \neq 0$. That $\gamma\left(-c_{\alpha}^{-1} d_{\alpha}\right) \neq 0$ and $\gamma\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right) \neq 0$ will be shown in the following proof.

Proof. Write $a, b, c, d$ for $a_{\alpha}, b_{\alpha}, c_{\alpha}, d_{\alpha}$; let $\sigma=r_{\Omega}(\alpha)$ with $\alpha \in \Omega_{\mathbf{A}} \cap P_{\mathbf{A}} C^{\theta}$; let $\lambda$ be the characteristic function of $\prod_{v \in \mathrm{~h}} L_{v}$. By (A2.22) and (A2.3c) we have

$$
\begin{aligned}
N & (\mathrm{il}(\alpha))^{1 / 2} h(\sigma, z)^{-1}=\left(\sigma \lambda_{\mathbf{A}}\right)(0 ; 0, z) \\
& =|\operatorname{det}(c)|_{\mathbf{A}}^{1 / 2} \int_{X_{\mathbf{A}}} \lambda_{\mathbf{A}}(y c ; 0, z) \mathbf{e}_{\mathbf{A}}\left(y c \cdot{ }^{t} d \cdot{ }^{t} y / 2\right) d y \\
& =|\operatorname{det}(c)|_{\mathbf{A}}^{-1 / 2} \int_{X_{\mathbf{A}}} \lambda_{\mathbf{A}}(x ; 0, z) \mathbf{e}_{\mathbf{A}}\left(x c^{-1} d \cdot{ }^{t} x / 2\right) d x \\
& =|\operatorname{det}(c)|_{\mathbf{A}}^{-1 / 2} \prod_{v \in \mathbf{a}} \int_{X_{v}} \varphi_{v}\left(x ; 0,\left(z+c^{-1} d\right)_{v}\right) d_{v} x \prod_{v \in \mathbf{h}} \int_{L_{v}} \mathbf{e}_{v}\left(x c^{-1} d \cdot{ }^{t} x / 2\right) d_{v} x .
\end{aligned}
$$

From (A2.10) we see that the integral over $X_{v}$ is equal to $\operatorname{det}\left(-i\left(z+c^{-1} d\right)_{v}\right)^{-1 / 2}$. The integral over $L_{v}$ is $N\left(\mathfrak{d}_{v}\right)^{-n / 2} \gamma_{v}\left(c^{-1} d\right)$. Our equality shows that $\gamma_{v}\left(c^{-1} d\right) \neq$ 0 . Since $\overline{\gamma_{v}(s)}=\gamma_{v}(-s)$, we obtain the first formula. We can take $\eta$ to be $\alpha$, since $\eta \in P_{\mathbf{A}} C^{\theta}$ as we noted it immediately after (A2.17). Then we find our assertion concerning $h(\eta, z)$. Now if $\rho=\operatorname{diag}\left[\delta 1_{n}, \delta^{-1} 1_{n}\right]$, then (8) of Theorem A2.4 together with (A2.3a) shows that ${ }^{\eta} \lambda=\prod_{v \in \mathbf{h}} r_{P}\left(\rho_{v}\right) \lambda_{v}=N(\mathfrak{d})^{-n / 2} \lambda^{\prime}$, where $\lambda^{\prime}(x)=\lambda(\delta x)$. Let $\sigma=r_{\Omega}\left(\alpha \eta^{-1}\right) \eta$ with $\alpha \in \Omega_{\mathbf{A}} \iota \cap P_{\mathbf{A}} C^{\theta}$. By (1) of Theorem A2.4 we have

$$
\begin{aligned}
\left(\sigma \lambda_{\mathbf{A}}\right)(x ; 0, z) & =\left(r_{\Omega}\left(\alpha \eta^{-1}\right) \eta \lambda_{\mathbf{A}}\right)(x ; 0, z) \\
& =N(\mathfrak{d})^{-n / 2} h(\eta, z)^{-1}\left(r_{\Omega}\left(\alpha \eta^{-1}\right) \lambda_{\mathbf{A}}^{\prime}\right)(x ; 0, \eta(z))
\end{aligned}
$$

Since $\alpha \eta^{-1}=\left[\begin{array}{ll}-b & a \\ -d & c\end{array}\right]$, a calculation similar to the above one shows that

$$
\begin{aligned}
\left(r_{\Omega}\left(\alpha \eta^{-1}\right) \lambda_{\mathbf{A}}^{\prime}\right)(0 ; 0, \eta(z))= & |\operatorname{det}(d)|_{\mathbf{A}}^{-1 / 2} N(\mathfrak{d})^{n / 2} \\
& \cdot \gamma\left(-\delta^{-2} d^{-1} c\right) \prod_{v \in \mathbf{a}} \operatorname{det}\left(i\left(z^{-1}+d^{-1} c\right)_{v}\right)^{-1 / 2}
\end{aligned}
$$

Since (A2.22) is nonzero, we see that $\gamma\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right) \neq 0$. Putting $z=r \mathbf{i}$ and taking the limit as $r \rightarrow 0$, we obtain the second formula.

A2.7. Proposition. Let $\psi^{*}$ be the quadratic ideal character of $F$ corresponding to the extension $F(\sqrt{-1}) / F$. Suppose $\alpha \in G \cap P_{\mathbf{A}} D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]$. Then $d_{\alpha}$ is invertible, $\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{D}}(\alpha)^{-1}$ is prime to 2, and

$$
h(\alpha, z)^{2}=\operatorname{sgn}\left(N_{F / \mathbf{Q}}\left(\operatorname{det}\left(d_{\alpha}\right)\right) \psi^{*}\left(\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{D}}(\alpha)^{-1}\right) j_{\alpha}(z)^{\mathbf{a}}\right.
$$

Proof. From Lemma 1.11 (2), (3) we obtain $\nu_{0}\left(\delta^{-1} d_{\alpha}^{-1} c_{\alpha}\right)=\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{d}}(\alpha)^{-1}$ and also the first two assertions of our proposition. Since $\alpha=\alpha \eta^{-1} \eta=r_{\Omega}\left(\alpha \eta^{-1}\right)$ $\cdot r_{\Omega}(\eta)$, Lemma A2.6 shows that

$$
\lim _{z \rightarrow 0} h(\alpha, z)^{2} /|h(\alpha, z)|^{2}=\omega\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right)^{2} .
$$

By Lemma A1.6 (2), $\prod_{v \nmid 2} \omega_{v}\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right)^{2}=\psi^{*}\left(\nu_{0}\left(\delta^{-1} d_{\alpha}^{-1} c_{\alpha}\right)\right)$. This combined with Theorem 2.4 (2) proves our last assertion, since $\lim _{z \rightarrow 0} j_{\alpha}(z)^{\mathbf{a}}=N_{F / \mathbf{Q}}\left(\operatorname{det}\left(d_{\alpha}\right)\right)$.

A2.8. Proposition. (1) If $\alpha \in \Gamma\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]$, then $\gamma_{v}\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right)=1$ for $v \mid 2$, $\operatorname{det}\left(d_{\alpha}\right) \mathfrak{g}$ is prime to 2 , and

$$
\begin{align*}
\lim _{z \rightarrow 0} h(\alpha, z) & =\left|N_{F / \mathbf{Q}}\left(\operatorname{det}\left(d_{\alpha}\right)\right)\right| \gamma\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right)  \tag{}\\
& =\sum_{x \in L / L d_{\alpha}} \mathbf{e}_{\mathbf{h}}\left(\delta^{-2} x d_{\alpha}^{-1} c_{\alpha} \cdot{ }^{t} x / 2\right) .
\end{align*}
$$

(2) Let $\beta=\xi \alpha \xi^{-1}$ with $\alpha$ as above and $\xi=\operatorname{diag}\left[1_{n}, s 1_{n}\right], 0 \ll s \in \mathfrak{g}$; suppose $\beta \in \Gamma\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]$. Let $\psi_{s}^{*}$ be the ideal character of $F$ corresponding to the extension $F(\sqrt{s}) / F$. Then $\operatorname{det}\left(d_{\alpha}\right) \mathfrak{g}$ is prime to $2 s$ and

$$
h(\beta, z)=h(\alpha, s z) \psi_{s}^{*}\left(\operatorname{det}\left(d_{\alpha}\right) \mathfrak{g}\right)
$$

Proof. Let $\alpha \in \Gamma\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]$; then $\operatorname{det}\left(d_{\alpha}\right) \mathfrak{g}$ is clearly prime to 2 and $\mathrm{il}_{\mathfrak{\jmath}}(\alpha)=$ $\mathfrak{g}$, so that $\nu_{0}\left(\delta^{-1} d_{\alpha}^{-1} c_{\alpha}\right)=\operatorname{det}\left(d_{\alpha}\right) \mathfrak{g}$ by Lemma 1.11 (2). Since $d_{\alpha}^{-1} c_{\alpha} \prec 2 \mathfrak{d}_{v}$ for $v \mid 2$, we see from (16.5) that $\gamma_{v}\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right)=1$ for $v \mid 2$. By Lemma A1.6 (3), $\left|\gamma_{v}\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right)\right|=\nu\left(\left(\delta^{-1} d_{\alpha}^{-1} c_{\alpha}\right)_{v}\right)^{-1 / 2}$ for $v \nmid 2$, and hence $\left|\gamma\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right)\right|=$ $N\left(\operatorname{det}\left(d_{\alpha}\right) \mathfrak{g}\right)^{-1 / 2}=\lim _{z \rightarrow 0}|h(\alpha, z)|^{-1}$. This combined with Lemma A2.6 proves our first equality of (*). By (16.5) we can easily express the quantity in question as a sum over $L / L d_{\alpha}$ as stated. Next, if $\beta$ is as in (2), then $c_{\beta}=s c_{\alpha}$ and $d_{\beta}=d_{\alpha}$, so that $h(\beta, z) / h(\alpha, s z)$ is a constant, which equals $\gamma\left(\delta^{-2} d_{\beta}^{-1} c_{\beta}\right) / \gamma\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right)$ by (1). If $v \nmid 2 \operatorname{det}\left(d_{\alpha}\right)$, then $d_{\alpha}^{-1}$ is $v$-integral, so that $\gamma_{v}\left(\delta^{-2} d_{\beta}^{-1} c_{\beta}\right)=\gamma_{v}\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right)=1$. Now $\operatorname{det}\left(d_{\alpha}\right) \mathfrak{g}$ is prime to $2 s$ since $b_{\beta} \prec 2 \mathfrak{d}^{-1}$ and $c_{\beta} \prec 2 s \mathfrak{d}$. If $v \mid \operatorname{det}\left(d_{\alpha}\right)$ and $v \nmid 2$, then by Lemma A1.6 (1) we have

$$
\gamma_{v}\left(\delta^{-2} d_{\beta}^{-1} c_{\beta}\right)=\gamma_{v}\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right)\left(\frac{s}{\nu_{0}\left(\delta^{-1} d_{\alpha}^{-1} c_{\alpha}\right)_{v}}\right)=\gamma_{v}\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right)\left(\frac{s}{\operatorname{det}\left(d_{\alpha}\right) \mathfrak{g}_{v}}\right)
$$

Therefore we obtain the formula of (2).
A2.9. Proof of Theorem 6.8. Observe that $\theta(u, z ; \ell)$ coincides with $\varphi_{F}(u, z ; \lambda)$ if $\lambda(x)=\ell\left({ }^{t} x\right)$. Then (1), (2), (3), and (5) of Theorem 6.8 are special cases of Theorem A2.4 (2), (4), (3), and Proposition A2.7, respectively. Given such $\lambda$ and $\ell$, put $D_{\lambda}=\left\{\left.\sigma \in C^{\theta}\right|^{\sigma} \ell=\ell\right\}$, where we write ${ }^{\sigma} \ell$ for ${ }^{\tau} \ell$ with any $\tau \in M_{\mathbf{A}}$ such that $\sigma=\operatorname{pr}(\tau)$. This is meaningful by Theorem A2.4 (5), and $D_{\lambda}$ is an open subgroup of $C^{\theta}$ by Theorem A2.4 (6); moreover $D_{\lambda}=C^{\theta}$ if $\ell$ is the characteristic function of $\prod_{v \in \mathbf{h}} L_{v}$. Now let the notation be as in (4) of Theorem 6.8. Then we can put $\alpha^{-1}=\gamma \sigma$ with $\gamma=\operatorname{diag}\left[{ }^{t} d, d^{-1}\right]$ and $\sigma \in D_{\lambda}$. Put $\ell_{1}={ }^{\alpha^{-1}} \ell$. By Theorem A2.4 (4) we have $\ell={ }^{\alpha} \ell_{1}$ and $\ell_{1}={ }^{\gamma}{ }^{\sigma} \ell={ }^{\gamma} \ell$, and so by Proposition A2.5,

$$
\theta(\alpha(u, z) ; \ell)=h_{\alpha}(z) \theta\left(u, z ;{ }^{\gamma} \ell\right)
$$

If $\ell$ is as in Theorem A2.4 (8), then by that assertion and (A2.3a) we have $\left({ }^{\gamma} \ell\right)(x)=$ $\ell\left(x \cdot{ }^{t} d\right)$, so that $\left({ }^{\gamma} \ell\right)\left({ }^{t} x\right)=\lambda(d x)$. Clearly this is valid for an arbitrary $\lambda \in \mathcal{S}\left(F_{\mathrm{h}}^{n}\right)$. Thus we obtain (4) of Theorem 6.8.

Proof of Theorem 6.9. The first part of (2) of Theorem 6.9 follows from Proposition A2.5. To prove (6.33), take $\alpha=\eta$ in Proposition A2.5 and substitute $\eta(u, z)$ for $(u, z)$. Then we obtain $\theta\left(-u, z ;{ }^{\eta} \ell\right)=h_{\eta}(\eta z) \theta(\eta(u, z) ; \ell)$. From (A2.24) we easily see that $h_{\eta}(\eta z)=h_{\eta}(z)^{-1}$, and hence $\theta(\eta(u, z) ; \ell)=h_{\eta}(z) \theta\left(-u, z ;{ }^{\eta} \ell\right)$. Now $\theta\left(-u, z ; \ell_{1}\right)=\theta\left(u, z ; \ell_{2}\right)$ with $\ell_{2}(x)=\ell_{1}(-x)$. Therefore $\theta(\eta(u, z) ; \ell)=h_{\eta}(z)$ $\cdot \theta\left(u, z ; \ell^{\prime}\right)$ with $\ell^{\prime}(x)=\left({ }^{\eta} \ell\right)(-x)$. By Theorem A2.4 (7) and (A2.3c) we obtain

$$
{ }^{\eta} \ell(x)=\left|D_{F}\right|^{-n / 2} \int_{\mathbf{F}_{\mathbf{h}}^{n}} \ell\left({ }^{t} y\right) \mathbf{e}_{\mathbf{h}}\left(-{ }^{t} x y\right) d y
$$

and hence we obtain $\lambda^{\prime}$ as given in (6.33).
To prove (1), we first consider the case $\alpha \in G$. Take any nonzero $\ell$ in $\mathcal{S}\left(X_{\mathrm{h}}\right)$ and define $\ell^{\prime}$ as in Proposition A2.5 with any choice of $\xi$ for a given $\alpha \in G$. (Take $\alpha$ as $\beta$ there.) We may assume that $r(z)=\prod_{v \in \mathbf{a}} g\left(\xi_{v}, z_{v}\right)$. By Theorem A2.4(6) we can find a congruence subgroup $\Gamma$ of $G$ such that ${ }^{\gamma} \ell=\ell$ and ${ }^{\gamma} \ell^{\prime}=\ell^{\prime}$ for every $\gamma \in \Gamma$. Let $\gamma \in \alpha^{-1} \Gamma \alpha \cap \Gamma$. By Proposition A2.5 we then have

$$
\begin{aligned}
& h\left(\alpha \gamma \alpha^{-1}, \alpha z\right) \theta\left(\alpha(u, z) ; \ell^{\prime}\right)=\theta\left(\alpha \gamma \alpha^{-1} \alpha(u, z) ; \ell^{\prime}\right)=r(\gamma z) \theta(\gamma(u, z) ; \ell) \\
& \quad=r(\gamma z) h(\gamma, z) \theta(u, z ; \ell)=r(\gamma z) h(\gamma, z) r(z)^{-1} \theta\left(\alpha(u, z) ; \ell^{\prime}\right) .
\end{aligned}
$$

Since $\theta\left(u, z ; \ell^{\prime}\right)$ is a nonzero function, we obtain (1) when $\alpha \in G$. To treat the general case, we observe that every element of $\widetilde{G}_{+}$is the product of an element of $G$ and an element of the form $\operatorname{diag}\left[1_{n}, s 1_{n}\right]$ with $s \in F^{\times}, \gg 0$. Therefore we can reduce our problem to the equality $h\left(\alpha \gamma \alpha^{-1}, \alpha z\right)=h_{\gamma}(z)$ for $\gamma$ in some congruence subgroup when $\alpha=\operatorname{diag}\left[1_{n}, s 1_{n}\right]$. We easily see that it is sufficient to prove it for $s \in \mathfrak{g}$. Taking $\alpha$ and $\gamma$ here to be $\xi$ and $\alpha$ of Proposition A2.8 (2), we btain $h\left(\alpha \gamma \alpha^{-1}, z\right)=h_{\gamma}(s z) \psi_{s}^{*}\left(\operatorname{det}\left(d_{\gamma}\right) \mathfrak{g}\right)$. Since $s \gg 0$, we have $\psi_{s}^{*}\left(\operatorname{det}\left(d_{\gamma}\right) \mathfrak{g}\right)=$ 1 if $\operatorname{det}\left(d_{\gamma}\right)-1 \in 4 s g$. Thus we can define the desired congruence subgroup by that congruence condition. This completes the proof of Theorem 6.9.

Let us now study the behavior of $h(\sigma, z)$ and the action of $\mathfrak{M}$ on $\mathcal{S}\left(V_{\mathbf{h}}\right)$ under the reflection $z \mapsto-z^{\rho}$, where $x^{\rho}$ is defined for $x \in\left(\mathbf{C}_{m}^{m}\right)^{\mathbf{a}}$ by $x_{v}^{\rho}=\bar{x}_{v}$ for each $v \in$ a. Observe that this reflection maps $\mathcal{H}$ onto itself. Putting $\alpha^{*}=E \alpha E^{-1}$ for $\alpha \in G_{\mathbf{A}}$ with $E=\operatorname{diag}\left[1_{n},-1_{n}\right]$, we see that $\left(C^{\theta}\right)^{*}=C^{\theta}, \alpha^{*}\left(-z^{\rho}\right)=-\alpha(z)^{\rho}$ for $z \in \mathcal{H}, \mu_{0}(\alpha, z)^{\rho}=\mu_{0}\left(\alpha^{*},-z^{\rho}\right)$, and $\overline{\mu(\alpha, z)}=\mu\left(\alpha^{*},-z^{\rho}\right)$.

A2.10. Proposition. There exists an automorphism of $M_{\mathbf{A}}$ which is written $\sigma \mapsto \sigma^{*}$, consistent with $\alpha \mapsto E \alpha E^{-1}$ for $\alpha \in G$, and determined by the relation $(\sigma f)^{*}=\sigma^{*} f^{*}$ for $f \in L^{2}\left(X_{\mathbf{A}}\right)$, where $f^{*}(x)=\overline{f(-x)}$. Moreover, $\operatorname{pr}(\sigma)^{*}=\operatorname{pr}\left(\sigma^{*}\right)$, $r_{P}(\alpha)^{*}=r_{P}\left(\alpha^{*}\right)$ for $\alpha \in P_{\mathbf{A}}, r_{\Omega}(\beta)^{*}=r_{\Omega}\left(\beta^{*}\right)$ for $\beta \in \Omega_{\mathbf{A}}, t^{*}=t^{-1}$ for $t \in \mathbf{T}$, $\mathfrak{M}^{*}=\mathfrak{M}, \overline{h(\sigma, z)}=h\left(\sigma^{*},-z^{\rho}\right)$ and $\sigma^{*}\left(\ell^{*}\right)=\left({ }^{\sigma} \ell\right)^{*}$ for every $\sigma \in \mathfrak{M}$ and $\ell \in \mathcal{S}\left(X_{\mathbf{h}}\right)$, where $\ell^{*}$ is defined by $\ell^{*}(x)=\overline{\ell(-x)}$.

Proof. From (A2.3a, c) we easily see that $\left[r_{P}(\alpha) f\right]^{*}=r_{P}\left(\alpha^{*}\right) f^{*}$ for $\alpha \in P_{\mathbf{A}}$ and $\left[r_{\Omega}(\beta) f\right]^{*}=r_{\Omega}\left(\beta^{*}\right) f^{*}$ for $\beta \in \Omega_{\mathbf{A}}$. From Lemma 7.5 we easily see that $M_{\mathbf{A}}$ is generated by $r_{P}\left(P_{\mathbf{A}}\right), r_{\Omega}\left(\Omega_{\mathbf{A}}\right)$, and $\mathbf{T}$, and hence these equalities prove our assertions except the last two, which follow from Theorem A2.4 (1) and (A2.22) combined with the relation $\left(\ell_{\mathbf{A}}\right)^{*}(x ; u, z)=\left(\ell^{*}\right)_{\mathbf{A}}\left(x ; \bar{u},-z^{\rho}\right)$ that can easily be verified.

A2.11. We conclude this section by investigating $h(\alpha, z)$ for $\alpha$ in the subgroup $S p(r, F) \times S p(s, F)$ of $S p(r+s, F)$. To emphasize the dimension, let us denote the symbols $G, P, X, \mathcal{H}, \mathcal{U}, M_{\mathbf{A}}$, and $\mathfrak{M}$ by $G^{(n)}, P^{(n)}, X^{(n)}, \mathcal{H}^{(n)}, \mathcal{U}^{(n)}, M_{\mathbf{A}}^{(n)}$, and $\mathfrak{M}^{(n)}$. Let $n=r+s$ with positive integers $r$ and $s$. For $\beta \in G_{\mathbf{A}}^{(r)}$ and $\gamma \in G_{\mathbf{A}}^{(s)}$ we define an element $\beta \times \gamma$ of $G_{\mathrm{A}}^{(n)}$ by (23.5).

Let us now study how this injection $(\beta, \gamma) \mapsto \beta \times \gamma$ of $G_{\mathbf{A}}^{(r)} \times G_{\mathbf{A}}^{(s)}$ into $G_{\mathbf{A}}^{(n)}$ can be extended to their metaplectic coverings. First, for $f \in L^{2}\left(X_{\mathbf{A}}^{(r)}\right)$ and $f^{\prime} \in L^{2}\left(X_{\mathbf{A}}^{(s)}\right)$ we define $f \otimes f^{\prime} \in L^{2}\left(X_{\mathbf{A}}^{(n)}\right)$ by $\left(f \otimes f^{\prime}\right)\left(x, x^{\prime}\right)=f(x) f^{\prime}\left(x^{\prime}\right)$ for $x \in X_{\mathbf{A}}^{(r)}$ and $x^{\prime} \in X_{\mathbf{A}}^{(s)}$. Given $\sigma \in M_{\mathbf{A}}^{(r)}$ and $\sigma^{\prime} \in M_{\mathbf{A}}^{(s)}$, we have a unique unitary operator $\left\langle\sigma, \sigma^{\prime}\right\rangle$ on $L^{2}\left(X_{\mathbf{A}}^{(n)}\right)$ such that $\left\langle\sigma, \sigma^{\prime}\right\rangle\left(f \otimes f^{\prime}\right)=\sigma f \otimes \sigma^{\prime} f^{\prime}$ for all such $f$ and $f^{\prime}$. Then, checking the formula of $[\mathrm{W} 64,(15)]$, we find that $\left\langle\sigma, \sigma^{\prime}\right\rangle \in M_{\mathbf{A}}^{(n)}$ and $\operatorname{pr}\left(\left\langle\sigma, \sigma^{\prime}\right\rangle\right)=\operatorname{pr}(\sigma) \times \operatorname{pr}\left(\sigma^{\prime}\right)$. (Notice that $\left(\sigma, \sigma^{\prime}\right) \mapsto\left\langle\sigma, \sigma^{\prime}\right\rangle$ is not injective.) For $z \in \mathcal{H}^{(r)}$ and $z^{\prime} \in \mathcal{H}^{(s)}$ define $\left[z, z^{\prime}\right] \in \mathcal{H}^{(n)}$ by $\left[z, z^{\prime}\right]_{v}=\operatorname{diag}\left[z_{v}, z_{v}^{\prime}\right]$; similarly for $u \in \mathcal{U}^{(r)}$ and $u^{\prime} \in \mathcal{U}^{(s)}$ define $\left[u, u^{\prime}\right] \in \mathcal{U}^{(n)}$ by $\left[u, u^{\prime}\right]_{v}={ }^{t}\left({ }^{t} u_{v},{ }^{t} u_{v}^{\prime}\right)$. Clearly

$$
\begin{equation*}
\varphi\left(x, x^{\prime} ;\left[u, u^{\prime}\right],\left[z, z^{\prime}\right]\right)=\varphi(x ; u, z) \varphi\left(x^{\prime} ; u^{\prime}, z^{\prime}\right) \tag{A2.25}
\end{equation*}
$$

Observe that $\left\langle\sigma, \sigma^{\prime}\right\rangle \in \mathfrak{M}^{(n)}$ if $\sigma \in \mathfrak{M}^{(r)}$ and $\sigma^{\prime} \in \mathfrak{M}^{(s)}$.
A2.12. Proposition. If $\sigma \in \mathfrak{M}^{(r)}$ and $\sigma^{\prime} \in \mathfrak{M}^{(s)}$, we have

$$
h\left(\left\langle\sigma, \sigma^{\prime}\right\rangle,\left[z, z^{\prime}\right]\right)=h(\sigma, z) h\left(\sigma^{\prime}, z^{\prime}\right) \quad \text { and } \quad\left\langle\sigma, \sigma^{\prime}\right\rangle\left(\ell \otimes \ell^{\prime}\right)=\left({ }^{\sigma} \ell\right) \otimes\left(\sigma^{\sigma^{\prime}} \ell^{\prime}\right)
$$

for $\ell \in \mathcal{S}\left(X_{\mathbf{h}}^{(r)}\right)$ and $\ell^{\prime} \in \mathcal{S}\left(X_{\mathbf{h}}^{(s)}\right)$, where $\left(\ell \otimes \ell^{\prime}\right)\left(x, x^{\prime}\right)=\ell(x) \ell^{\prime}\left(x^{\prime}\right)$. Moreover, if $\alpha \in G^{(r)}$ and $\alpha^{\prime} \in G^{(s)}$, then $\alpha \times \alpha^{\prime}$ as an element of $M_{\mathrm{A}}^{(n)}$ coincides with $\left\langle\alpha, \alpha^{\prime}\right\rangle$, and hence

$$
h\left(\alpha \times \alpha^{\prime},\left[z, z^{\prime}\right]\right)=h(\alpha, z) h\left(\alpha^{\prime}, z^{\prime}\right)
$$

if $\alpha \in G^{(r)} \cap \mathfrak{M}^{(r)}$ and $\alpha^{\prime} \in G^{(s)} \cap \mathfrak{M}^{(s)}$.
Proof. Formula (A2.25) together with (A2.22) proves the first equality, which together with Theorem A2.4 (1) proves the second one. Let $\alpha \in G^{(r)}$ and $\alpha^{\prime} \in G^{(s)}$. If $\alpha \in P^{(r)}$ and $\alpha^{\prime} \in P^{(s)}$, formula (A2.3a) gives $\alpha \times \alpha^{\prime}=\left\langle\alpha, \alpha^{\prime}\right\rangle$. Since $G$ is generated by $P$ and $\eta$, our proof will be complete if we can show $\eta_{r} \times 1=\left\langle\eta_{r}, 1\right\rangle$ and $1 \times \eta_{s}=\left\langle 1, \eta_{s}\right\rangle$. Now, from Lemma A2.6 we easily see that

$$
\begin{equation*}
h\left(\alpha \times 1,\left[z, z^{\prime}\right]\right)=h(\alpha, z) \quad \text { if } \quad \alpha \in G^{(r)} \quad \text { and } \quad \operatorname{det}\left(d_{\alpha}\right) \neq 0 . \tag{A2.26}
\end{equation*}
$$

Observe that $\eta_{r}=\beta \gamma$ with $\beta=\left[\begin{array}{cc}0 & -1_{r} \\ 1_{r} & 2 \cdot 1_{r}\end{array}\right]$ and $\gamma=\left[\begin{array}{cc}1_{r} & -2 \cdot 1_{r} \\ 0 & 1_{r}\end{array}\right]$. Clearly $\operatorname{det}\left(d_{\beta}\right) \operatorname{det}\left(d_{\gamma}\right) \neq 0$, and hence the equality of (A2.26) is true with $\beta$ and $\gamma$ in place of $\alpha$. Since $\gamma \in C^{\theta}$, Theorem A2.4 (4) shows that
$h\left(\eta_{r}, z\right)=h(\beta, \gamma z) h(\gamma, z)=h\left(\beta \times 1,\left[\gamma z, z^{\prime}\right]\right) h\left(\gamma \times 1,\left[z, z^{\prime}\right]\right)=h\left(\eta_{r} \times 1,\left[z, z^{\prime}\right]\right)$.
On the other hand we already know that $h\left(\left\langle\eta_{r}, 1\right\rangle,\left[z, z^{\prime}\right]\right)=h\left(\eta_{r}, z\right)$. Observe that if $h(\sigma, w)=h(\tau, w)$ for some $w$ and $\operatorname{pr}(\sigma)=\operatorname{pr}(\tau)$, then $\sigma=\tau$. Therefore we have $\left\langle\eta_{r}, 1\right\rangle=\eta_{r} \times 1$. Similarly $\left\langle 1, \eta_{s}\right\rangle=1 \times \eta_{s}$, which completes the proof.

A2.13. Proof of Proposition 16.9. We consider only half-integral $k$; the case of integral $k$ can be handled by the same methods (see Remark 16.12). We identify any element of $P_{\mathbf{A}}$ with its image under $r_{P}$, as we did in $\S 16.5$. We start with an obvious equality

$$
\begin{equation*}
E_{\mathbf{A}}^{*}(x, s)=\chi(\delta)^{-n} \sum_{\alpha \in A} \mu(\alpha x \widetilde{\zeta}) \varepsilon(\alpha x \widetilde{\zeta})^{-s}, \quad A=P \backslash G \tag{i}
\end{equation*}
$$

Let $\alpha \in G$ and $\operatorname{pr}(x) \in P_{\mathbf{A}} G_{\mathbf{a}}$; suppose $\mu(\alpha x \widetilde{\zeta}) \neq 0$. Then $\alpha x \widetilde{\zeta} \in P_{\mathbf{A}} \widetilde{D}$, so that $\alpha \in P_{\mathbf{A}} D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right] \zeta^{-1} P_{\mathbf{A}}$. Now (1.18) shows that $\operatorname{det}\left(c_{y}\right)_{v} \neq 0$ if $y \in D\left[\mathfrak{b}^{-1}, \mathfrak{b} \mathfrak{c}\right] \zeta^{-1}$ and $v \mid c$. Therefore we see that $\operatorname{det}\left(c_{\alpha}\right) \neq 0$, and so $\alpha \in P \eta P$ by [S97, Lemma 2.12 (2)], where $\eta$ is $\eta_{n}$ of (1.8). Thus we can take $P \backslash P \eta P$ in place of $A$ in (i). By [S97, Lemma 18.8 (2)] we can take $\eta R$ as $P \backslash P \eta P$, where $R$ is given by (16.11). Thus

$$
\begin{equation*}
E^{*}(x, s)=\chi(\delta)^{-n} \sum_{\alpha \in \eta R} \mu(\alpha x \widetilde{\zeta}) \varepsilon(\alpha x \widetilde{\zeta})^{-s} \text { if } \operatorname{pr}(x) \in P_{\mathbf{A}} G_{\mathbf{a}} \tag{ii}
\end{equation*}
$$

For $\sigma \in S_{\mathbf{A}}$ put $\tau(\sigma)=\left[\begin{array}{cc}1_{n} & \sigma \\ 0 & 1_{n}\end{array}\right]$. Clearly $\tau(S)=R$. Putting $\xi=\operatorname{diag}[q, \hat{q}]$, from (16.42) we obtain

$$
\begin{equation*}
c(h, q, s)=\chi(\delta)^{-n} \int_{S_{\mathbf{A}} / S} E^{*}(\tau(\sigma) \xi, s) \mathbf{e}_{\mathbf{A}}^{n}(-h \sigma) d \sigma \tag{iii}
\end{equation*}
$$

where we take the measure of $S_{\mathbf{A}} / S$ to be 1 . To simplify our notation, put $g(x)=$ $\mu(x) \varepsilon(x)^{-s}$. Putting $x=\tau(\sigma) \xi$ and $\alpha=\eta \tau(a)$ with $a \in S$ in (ii), we obtain

$$
E^{*}(\tau(\sigma) \xi, s)=\chi(\delta)^{-n} \sum_{a \in S} g(\eta \tau(a+\sigma) \xi \widetilde{\zeta})
$$

Substituting this into (iii), we find that

$$
\begin{equation*}
c(h, q, s)=\chi(\delta)^{-n} \int_{S_{\mathbf{A}}} g(\eta \tau(\sigma) \xi \widetilde{\zeta}) \mathbf{e}_{\mathbf{A}}^{n}(-h \sigma) d \sigma \tag{iv}
\end{equation*}
$$

Since $\tau(b) \xi=\xi \tau\left(q^{-1} b \widehat{q}\right)$, from (16.41) we see that $E^{*}(\tau(\sigma+b) \xi)=E^{*}(\tau(\sigma) \xi)$ if $b \in S_{\mathbf{h}}$ and $q^{-1} b \widehat{q} \prec \mathfrak{d}^{-2} b c$. Then (iii) shows that $c(h, q, s) \neq 0$ only if $\left({ }^{t} q h q\right)_{v} \in$ $\left(\mathfrak{d} \mathfrak{b}^{-1} \mathfrak{c}^{-1}\right)_{v} \widetilde{S}_{v}$ for every $v \in \mathbf{h}$, which is the first statement of Proposition 16.9. Consequently $\alpha_{c}^{\lambda}\left(\varepsilon_{b}^{-1} \cdot{ }^{t} q h q, 2 s, \chi\right)$ is meaningful.

Our next task is to determine the value of $g(x \widetilde{\zeta})$ for $x=\eta \tau(\sigma) \xi$. Putting $y=x \zeta$, we have

$$
y_{\mathbf{h}}=\left[\begin{array}{cc}
-\delta \widehat{q} & 0 \\
\delta \sigma \widehat{q} & -\delta^{-1} q
\end{array}\right]_{\mathbf{h}}, \quad y_{\mathbf{a}}=x_{\mathbf{a}}=\left[\begin{array}{cc}
0 & -\widehat{q} \\
q & \sigma \widehat{q}
\end{array}\right]_{\mathbf{a}}, \quad\left(d_{y}^{-1} c_{y}\right)_{\mathbf{h}}=-\delta^{2}\left(q^{-1} \sigma \widehat{q}\right)_{\mathbf{h}}
$$

We assume $\operatorname{det}\left(q_{v}\right)>0$ for every $v \in$ a. Since $r_{\Omega}(\eta)=\eta$ and we are identifying $\tau(\sigma) \xi$ with $r_{P}(\tau(\sigma) \xi)$, we have, by (A2.3b), $x=r_{\Omega}(\eta) r_{P}(\tau(\sigma) \xi)=r_{\Omega}(\eta \tau(\sigma) \xi)=$ $r_{\Omega}(x)$. Now $j^{k}(x \widetilde{\zeta}, \mathbf{i})=j_{x}^{k}(\mathbf{i})$ by $(16.30)$, and $j_{x}^{k}(\mathbf{i})=\lambda \operatorname{det}(q i+\sigma \widehat{q})_{\mathbf{a}}^{k}$ with $\lambda \in \mathbf{T}$. (The branch of $\operatorname{det}(z)^{k}$ was chosen in §16.8.) Then from Lemma A2.6 we can easily derive that $\lambda \mathbf{e}(n[F: \mathbf{Q}] / 8)=\omega\left(-q^{-1} \sigma \widehat{q}\right)$. Since $\omega(-s)=\omega(s)^{-1}$, we have $j_{x}^{k}(\mathbf{i})^{-1}=C \omega\left(q^{-1} \sigma \widehat{q}\right) \operatorname{det}(q i+\sigma \widehat{q})_{\mathbf{a}}^{-k}$ with $C=\mathbf{e}(n[F: \mathbf{Q}] / 8)$.

Now $g(x \widetilde{\zeta}) \neq 0$ only if $y \in P_{\mathbf{A}} D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]$, which is so if and only if $\left(\delta^{2} q^{-1} \sigma \widehat{q}\right)_{v} \prec$ $\mathfrak{b}_{v} \mathfrak{c}_{v}$ for every $v \mid \mathfrak{c}$, by virtue of the characterization of $P_{\mathbf{A}} D$ in Lemma 1.9. Assuming the last condition, let $y_{\mathbf{h}}=p w$ with $p \in P_{\mathbf{h}}$ and $w \in D\left[\mathfrak{b}^{-1}, \mathfrak{b} c\right]$; take $b \in F_{\mathbf{h}}^{\times}$ so that $b \mathfrak{g}=\mathfrak{b}$. By Lemma 1.11 (2),

$$
\operatorname{det}\left(d_{y} d_{p}^{-1}\right) \mathfrak{g}=\operatorname{det}\left(d_{y}\right) \mathrm{il}_{\mathfrak{b}}(y)^{-1}=\nu_{0}\left(b^{-1} d_{y}^{-1} c_{y}\right)=\nu_{0}\left(b^{-1} \delta^{2} q^{-1} \sigma \widehat{q}\right)
$$

Since $d_{y}=d_{p} d_{w}$, we obtain

$$
\begin{align*}
\chi_{\mathbf{h}}\left(\operatorname{det}\left(d_{p}\right)\right)^{-1} \chi_{\mathfrak{c}}\left(\operatorname{det}\left(d_{w}\right)\right)^{-1} & =\chi_{\mathbf{h}}\left(\operatorname{det}\left(d_{y}\right)\right)^{-1}\left(\chi_{\mathbf{h}} / \chi_{\mathbf{c}}\right)\left(\operatorname{det}\left(d_{p}^{-1} d_{y}\right)\right)  \tag{v}\\
& =\chi_{\mathbf{h}}\left(\operatorname{det}\left(-\delta q^{-1}\right)\right) \chi^{*}\left(\nu_{0}\left(b^{-1} \delta^{2} q^{-1} \sigma \widehat{q}\right)\right)
\end{align*}
$$

Here notice that $\nu_{0}\left(b^{-1} d_{y}^{-1} c_{y}\right)$ is prime to $\mathfrak{c}$ by Lemma 1.11 (3). By (16.23a), $\varepsilon\left(y_{\mathbf{h}}\right)=N\left(\mathrm{il}_{\mathbf{b}}(y)\right)^{-2}=\left|\delta^{-n} \operatorname{det}\left(q_{\mathbf{h}}\right)\right|_{\mathbf{A}}^{2} \nu\left(b^{-1} \delta^{2} q^{-1} \sigma \widehat{q}\right)^{2}$ and $\varepsilon\left(y_{\mathbf{a}}\right)=\varepsilon\left(x_{\mathbf{a}}\right)=\left|j_{x}(\mathbf{i})\right|^{2 \mathbf{a}}$. Temporarily denote the quantity of (v) by $\mu\left(y_{\mathrm{h}}\right)$. Combining all these and taking $\varepsilon_{b}$ as in Proposition 16.9, we obtain

$$
\begin{aligned}
\chi(\delta)^{-n} g(x \widetilde{\zeta})= & \chi(\delta)^{-n} \mu\left(y_{\mathbf{h}}\right) \varepsilon\left(y_{\mathbf{h}}\right)^{-s} \varepsilon\left(x_{\mathbf{a}}\right)^{-s} j_{x}^{k}(\mathbf{i})^{-1}\left|j_{x}(\mathbf{i})\right|^{k-i \kappa} \\
= & C\left|\delta^{n} \operatorname{det}\left(q_{\mathbf{h}}^{-1}\right)\right|_{\mathbf{A}}^{2 s} \chi_{\mathbf{h}}\left(\operatorname{det}\left(-q^{-1}\right)\right) \chi^{*}\left(\nu_{0}\left(\delta \varepsilon_{b} q^{-1} \sigma \widehat{q}\right)\right) \nu\left(\delta \varepsilon_{b} q^{-1} \sigma \widehat{q}\right)^{-2 s} \\
& \cdot \omega\left(q^{-1} \sigma \widetilde{q}\right) \operatorname{det}(q i+\sigma \widehat{q})_{\mathbf{a}}^{-k}\left|\operatorname{det}(q i+\sigma \widehat{q})_{\mathbf{a}}\right|^{k-i \kappa-2 s \mathbf{a}} .
\end{aligned}
$$

For a fixed $q$ the last product can be written in the form

$$
C\left|\delta^{n} \operatorname{det}(q)\right|_{\mathbf{A}}^{2 s} \chi_{\mathbf{h}}\left(\operatorname{det}\left(-q^{-1}\right)\right) \prod_{v \in \mathbf{v}} f_{v}\left(\sigma_{v}\right)
$$

with functions $f_{v}$ on $S_{v}$ which we choose in the manner obvious from the above expression. Then from (iv) we obtain

$$
c(h, q, s)=C\left|\delta^{n} \operatorname{det}(q)\right|_{\mathbf{A}}^{2 s} \chi_{\mathbf{h}}\left(\operatorname{det}\left(-q^{-1}\right)\right) c(S) \prod_{v \in \mathbf{v}} \int_{S_{v}} f_{v}\left(\sigma_{v}\right) \mathbf{e}_{v}^{n}\left(-h \sigma_{v}\right) d \sigma_{v}
$$

where $c(S)$ is the factor determined by $d \sigma=c(S) \prod_{v \in \mathrm{v}} d \sigma_{v}$, whose value is given in [S97, (18.9.3)]. Since $\left(\delta^{2} q^{-1} \sigma \widehat{q}\right)_{v} \prec \mathfrak{b}_{v} \mathfrak{c}_{v} \subset 2 \mathfrak{d}_{v}$ for every $v \mid \mathfrak{c}$, we have $\nu_{0}\left(\delta \varepsilon_{b} q^{-1} \sigma \widehat{q}\right)_{v}$ $=\mathfrak{g}_{v}$ and $\left(q^{-1} \sigma \widehat{q}\right)_{v} \prec 2 \mathfrak{d}_{v}^{-1}$ for $v \mid \mathfrak{c}$. Therefore, from (16.5) and (16.6) we obtain $\omega\left(\left(q^{-1} \sigma \widehat{q}\right)_{v}\right)=1$ for $v \mid \mathbf{c}$. Thus, for $v \mid c$, we can take $f_{v}$ to be the characteristic function of $q_{v} S\left(\mathfrak{d}^{-2} \mathfrak{b c}\right)_{v} \cdot{ }^{t} q_{v}$. Put $\mu_{v}=N\left(\mathfrak{b}_{v} \mathfrak{c}_{v}\right)^{-n(n+1) / 2}\left|\operatorname{det}\left(q_{v}\right)\right|^{n+1}$. Since $\left({ }^{t} q h q\right)_{v} \in\left(\mathfrak{d b}^{-1} \mathfrak{c}^{-1}\right)_{v} \widetilde{S}_{v}$, the integral over $S_{v}$ for $v \mid \mathfrak{c}$ is the measure of the set $q_{v} S\left(\mathfrak{d}^{-2} \mathfrak{b c}\right)_{v} \cdot{ }^{t} q_{v}$, which equals $\mu_{v} N\left(\mathfrak{d}_{v}\right)^{n(n+1)}$.

For $v \in \mathbf{h}, v \nmid c$, taking the variable change $\sigma_{v} \mapsto b_{v} q_{v} \sigma_{v} \cdot{ }^{t} q_{v}$, we find that

$$
\int_{S_{v}} f_{v}\left(\sigma_{v}\right) \mathbf{e}_{v}^{n}\left(-h \sigma_{v}\right) d \sigma_{v}=\mu_{v} \cdot \int_{S_{v}} \chi^{*}\left(\nu_{0}\left(\sigma_{v}\right)\right) \omega\left(\delta_{v}^{-1} \sigma_{v}\right) \nu\left(\sigma_{v}\right)^{-2 s} \mathbf{e}_{v}^{n}\left(-\tau \sigma_{v}\right) d \sigma_{v}
$$

where $\tau=\left(\varepsilon_{b}^{-1} \cdot{ }^{t} q h q\right)_{v}$. The last integral is clearly the $v$-factor of $\alpha_{\mathrm{c}}^{1}(\tau, 2 s, \chi)$.
As for $v \in \mathbf{a}$, the integral over $S_{v}$ produces $\xi\left(y_{v}, h_{v} ; \ldots\right)$. We shall not go into details here, since such a calculation is practically the same as what was done in [S97, $\S 18.11]$ for integral $k$. Taking the product of all these factors, we can complete the proof. As we said at the beginning, the case of integral $k$ can be proved in a similar and much simpler way.

A2.14. Proof of Theorem 17.7 (vi). Suppose $n=1$ and $\mathfrak{c}=\mathfrak{g}$; then $k \in \mathbf{Z a}^{\mathbf{a}}$ and $G=P \cup P \eta P$ by [S97, Lemma 2.12 (1)]. Since $P \backslash P \eta P$ can be given by $\eta R$ by [S97, Lemma 18.8 (2)], $P \backslash G$ can be given by $\{1\} \cup \eta R$. Thus, by (16.33), for $\xi \in G_{\mathrm{a}}$ we have

$$
\begin{gathered}
E_{\mathbf{A}}^{*}(\xi, s)=\chi(\delta)^{-1} \mu(\xi \zeta) \varepsilon(\xi \zeta)^{-s}+E^{\prime}(\xi, s) \\
\text { with } \quad E^{\prime}(\xi, s)=\chi(\delta)^{-1} \sum_{\alpha \in \eta R} \mu(\alpha \xi \zeta) \varepsilon(\alpha \xi \zeta)^{-s}
\end{gathered}
$$

Given $z=x+i y \in \mathfrak{H}_{1}^{\text {a }}$, take $\xi=\tau(x) \operatorname{diag}\left[y^{1 / 2}, y^{-1 / 2}\right]$. We now apply our computation of $\S A 2.13$ (for $k \in \mathbf{Z}^{\mathbf{a}}$, or that of $[\mathrm{S} 97, \S 18.11]$ ) to $E^{\prime}(\xi, s)$ to obtain

$$
E^{\prime}(\xi, s) j_{\xi}^{k}(\mathbf{i})=y^{-k / 2} \sum_{h \in F} c\left(h, y^{1 / 2}, s\right) \mathbf{e}_{\mathbf{a}}(h x)
$$

with quantities $c\left(h, y^{1 / 2}, s\right)$, to which the formula of Proposition 16.9 is applicable; we simply take $\mathfrak{c}=\mathfrak{g}$ in that formula. Observing that $\zeta$ of (16.28) belongs to
$\operatorname{diag}\left[\varepsilon_{b}^{-1}, \varepsilon_{b}\right] D\left[\mathfrak{b}^{-1}, \mathfrak{b}\right]$ with $\varepsilon_{b}$ of Proposition 16.9 , we see that $\chi(\delta)^{-1} \mu(\xi \zeta) \varepsilon(\xi \zeta)^{-s}$ $=\chi^{*}\left(\mathfrak{b} \mathfrak{d}^{-2}\right) N\left(\mathfrak{b}^{-1} \mathfrak{d}\right)^{2 s} y^{s \mathbf{a}}$. Therefore

$$
\begin{align*}
L(2 s, \chi) E^{*}(z, s)= & \chi^{*}\left(\mathfrak{b} \boldsymbol{d}^{-2}\right) N\left(\mathfrak{b}^{-1} \mathfrak{d}\right)^{2 s} L(2 s, \chi) y^{s \mathbf{a}-k / 2}  \tag{}\\
& +L(2 s, \chi) y^{-k / 2} \sum_{h \in F} c\left(h, y^{1 / 2}, s\right) \mathbf{e}_{\mathbf{a}}(h x)
\end{align*}
$$

Our problem is the nature of $D(z, 0 ; 2 \mathbf{a}, \chi, \mathfrak{g})$ for $n=1$. If $\chi \neq 1$, then (v) of our theorem says that it belongs to $\pi^{d} \mathcal{M}_{2 \mathrm{a}}\left(\mathbf{Q}_{\mathrm{ab}}\right)$. Suppose $\chi=1$; then $L(2 s, \chi)=$ $\zeta_{F}(2 s)$ and the analysis of the Fourier coefficients at $s=0$ in the setting of (v) is applicable to the second term of $\left(^{*}\right)$. Thus $\left(^{*}\right)$ with $k=2 \mathbf{a}$ at $s=0$ gives

$$
\zeta_{F}(0) y^{-\mathbf{a}}+\pi^{d} \sum_{h \in F} a_{h} \mathbf{e}_{\mathbf{a}}(h z)
$$

with $a_{h} \in \mathbf{Q}_{\mathrm{ab}}$. By Lemma $17.5(3), \zeta_{F}(0)=0$ if $F \neq \mathbf{Q}$. It is well-known that $\zeta(0)=-1 / 2$. Thus, for $k=2 \mathbf{a}$ and $s=0,\left(^{*}\right)$ produces an element of $\pi^{d} \mathcal{M}_{2 \mathbf{a}}\left(\mathbf{Q}_{\mathbf{a b}}\right)$ or $\pi \mathcal{N}_{2}^{1}\left(\mathbf{Q}_{\mathrm{ab}}\right)$ according as $F \neq \mathbf{Q}$ or $F=\mathbf{Q}$. This result is for $E^{*}$. Transforming $E^{*}$ back to $E$ by $\zeta_{0}$ as in (16.35), we can complete the proof.

## A3. Transformation formulas of general theta series

A3.1. For $X=\left(x_{i j}\right) \in \mathbf{C}_{q}^{q}$ and $Y \in \mathbf{C}_{n}^{n}$ we define $X \otimes Y \in \mathbf{C}_{n q}^{n q}$ by

$$
X \otimes Y=\left[\begin{array}{ccc}
x_{11} Y & \cdots & x_{1 q} Y \\
\cdots & \cdots & \cdots \\
x_{q 1} Y & \cdots & x_{q q} Y
\end{array}\right]
$$

Let $V$ be a $q$-dimensional vector space over $F$ and $S: V \times V \rightarrow F$ a nondegenerate $F$-bilinear symmetric form. For each $v \in \mathbf{v}$ we have an $F_{v}$-bilinear symmetric form $S_{v}: V \times V \rightarrow F_{v}$. For each $v \in \mathbf{a}$ put $I_{v}=\operatorname{diag}\left[1_{r_{v}},-1_{s_{v}}\right]$ with the signature $\left(r_{v}, s_{v}\right.$ ) of $S_{v}$; we also take and fix an $F_{v}$-linear bijection $A_{v}: V_{v} \rightarrow$ $\left(F^{q}\right)_{v}$ so that $S_{v}(x, y)={ }^{t}\left(A_{v} x\right) I_{v}\left(A_{v} y\right)$ and put $T_{v}(x, y)={ }^{t}\left(A_{v} x\right)\left(A_{v} y\right)$. For $p=$ $\left(p_{1}, \ldots, p_{n}\right) \in V_{v}^{n}$ with $p_{i} \in V_{v}$ we define elements $S_{v}[p]$ and $T_{v}[p]$ of $\left(F_{v}\right)_{n}^{n}$ by $S_{v}[p]=\left(S_{v}\left(p_{i}, p_{j}\right)\right)_{i, j=1}^{n}, \quad T_{v}[p]=\left(T_{v}\left(p_{i}, p_{j}\right)\right)_{i, j=1}^{n}$. Notice that $S_{v}[p]$ is meaningful also for $v \in \mathbf{h}$.

A3.2. Let us again emphasize the dimension as we did in Section A2, by using the symbols $G^{(n)}, \mathcal{H}^{(n)}, \mathcal{U}^{(n)}$; in addition we use $D^{(n)}[\mathfrak{x}, \mathfrak{y}]$ and $\theta^{(n)}$ for $D[\mathfrak{x}, \mathfrak{y}]$ and $\theta$. We now define a theta function $g(u, z ; \lambda)$ for $z \in \mathcal{H}^{(n)}, u \in \mathcal{U}^{(n q)}$, and $\lambda \in \mathcal{S}\left(V_{\mathbf{h}}^{n}\right)$ by

$$
\begin{equation*}
g(u, z ; \lambda)=\sum_{\xi \in V^{n}} \lambda\left(\xi_{\mathbf{h}}\right) \Phi(\xi ; u, z) \tag{A3.0a}
\end{equation*}
$$

$$
\begin{align*}
& \begin{array}{l}
\text { (A3.0b) } \\
\Phi(p ; u, z)=\prod_{v \in \mathbf{a}} \Phi_{v}\left(p_{v} ; u_{v}, z_{v}\right) \quad\left(p \in V_{\mathbf{A}}^{n}\right) \\
\Phi_{v}(p ; u, x+i y)=\mathbf{e}\left({ }^{t} u\left(1_{q} \otimes 4 i y\right)^{-1} u+2^{-1} \operatorname{tr}\left(x S_{v}[p]+i y T_{v}[p]\right)+\operatorname{tr}\left(u^{\prime}\left(A_{v} p\right)\right)\right)
\end{array} \tag{A3.0b}
\end{align*}
$$

$$
\text { for } p \in V_{v}^{n}, x+i y \in \mathfrak{H}_{n}, \text { and } u \in \mathbf{C}^{n q}
$$

where $A_{v} p=\left(A_{v} p_{1} \ldots A_{v} p_{n}\right) \in\left(F_{v}\right)_{n}^{q}$ for $p=\left(p_{1}, \ldots, p_{n}\right) \in V_{v}^{n}, u^{\prime}=\left(u_{1} \ldots u_{q}\right)$ $\in \mathbf{C}_{q}^{n}$ for ${ }^{t} u=\left({ }^{t} u_{1} \ldots{ }^{t} u_{q}\right) \in \mathbf{C}_{n q}^{1}$ with $u_{i} \in \mathbf{C}^{n}$. If $q=1, V=F, S(x, x)=x^{2}$, and $A_{v}=1$ for all $v$, then we see that $g(u, z ; \lambda)$ coincides with $\theta^{(n)}(u, z ; \lambda)$. In
the general case $g$ can be obtained as a "pullback" of $\theta^{(n q)}$ as will be shown below. It should be noted that $g(u, z ; \lambda)=0$ for every $(u, z)$ only if $\lambda=0$.

We let $G_{\mathbf{A}}$ act on $\mathcal{H}^{(n)} \times \mathcal{U}^{(n q)}$ by $\alpha(u, z)=(w, \alpha z)$ for $\alpha \in G_{\mathbf{A}}$ with

$$
\begin{equation*}
\left.w_{v}=\operatorname{diag}\left[1_{r_{v}} \otimes^{t} \mu(\alpha, z)_{v}^{-1}, 1_{s_{v}} \otimes^{t} \overline{\mu(\alpha, z}\right)_{v}^{-1}\right] u_{v} \tag{A3.1}
\end{equation*}
$$

We now put $\mathfrak{M}_{q}=\mathfrak{M}$ if $q$ is odd and $\mathfrak{M}_{q}=G_{\mathbf{A}}$ if $q$ is even, and define a factor of automorphy $J^{S}(\alpha, z)$ for $\alpha \in \mathfrak{M}_{q}$ by

$$
J^{S}(\alpha, z)= \begin{cases}\prod_{v \in \mathbf{a}} j_{\alpha}(z)_{v}^{(q / 2)-s_{v}}\left|j_{\alpha}(z)_{v}\right|^{s_{v}} & \text { if } q \text { is even }  \tag{A3.2a}\\ h(\alpha, z)^{q} \prod_{v \in \mathbf{a}} j_{\alpha}(z)_{v}^{-s_{v}}\left|j_{\alpha}(z)_{v}\right|^{s_{v}} & \text { if } q \text { is odd. }\end{cases}
$$

If $q$ is even, $J^{S}$ is a factor of automorphy; if $q$ is odd, however, Theorem A2.4 (4) implies the following weaker property:

$$
\begin{align*}
& J^{S}(\alpha \beta \gamma, z)=J^{S}(\alpha, z) J^{S}(\beta, \gamma z) J^{S}(\gamma, z)  \tag{A3.2b}\\
& \text { if } \operatorname{pr}(\alpha) \in P_{\mathbf{A}}, \beta \in \mathfrak{M}, \text { and } \operatorname{pr}(\gamma) \in C^{\theta} .
\end{align*}
$$

A3.3. Theorem. Let $\chi$ be the Hecke character of $F$ corresponding to the extension $F\left(\operatorname{det}(S)^{1 / 2}\right) / F$ or $F\left((-1)^{q / 4} \operatorname{det}(S)^{1 / 2}\right)$ according as $q$ is odd or even; let pr denote the identity map of $G_{\mathbf{A}}$ onto itself if $q$ is even. Then every $\sigma \in \mathfrak{M}_{q}$ gives a C-linear automorphism $\lambda \mapsto{ }^{\sigma} \lambda$ of $\mathcal{S}\left(V_{\mathrm{h}}^{n}\right)$ with the following properties:
(0) $J^{S}(\alpha, z)^{-1} g\left(\alpha(u, z) ;{ }^{\alpha} \lambda\right)=g(u, z ; \lambda)$ if $\alpha \in G \cap \mathfrak{M}_{q}$.
(1) The map $\lambda \mapsto{ }^{\sigma} \lambda$ does not depend on $\left\{A_{v}\right\}_{v \in \mathbf{a}}$ (though it depends on $S$ ).
(2) ${ }^{(\sigma \tau)} \lambda={ }^{\sigma}\left({ }^{\tau} \lambda\right)$ for every $\sigma$ and $\tau$ in $G_{\mathbf{A}}$ if $q$ is even; ${ }^{(\rho \sigma \tau)} \lambda={ }^{\rho}\left({ }^{\sigma}\left({ }^{\tau} \lambda\right)\right)$ whenever $\operatorname{pr}(\rho) \in P_{\mathbf{A}}$ and $\operatorname{pr}(\tau) \in C^{\theta}$ if $q$ is odd.
(3) ${ }^{\sigma} \lambda$ depends only on $\lambda$ and $\operatorname{pr}(\sigma)_{\mathbf{h}}$.
(4) $\operatorname{pr}\left(\left\{\left.\sigma \in \mathfrak{M}_{q}\right|^{\sigma} \lambda=\lambda\right\}\right)$ contains an open subgroup of $G_{\mathbf{A}}$ for every $\lambda \in \mathcal{S}\left(V_{\mathbf{h}}^{n}\right)$.
(5) For $\operatorname{pr}(\sigma)=\tau \in P_{\mathrm{h}}$ we have

$$
\left({ }^{\sigma} \lambda\right)(x)=\left|\operatorname{det}\left(a_{\tau}\right)_{\mathbf{h}}\right|_{\mathbf{A}}^{q / 2} \chi_{\mathbf{h}}\left(\operatorname{det}\left(a_{\tau}\right)\right) \mathbf{e}_{\mathbf{h}}\left(\operatorname{tr}\left(S[x] a_{\tau} \cdot{ }^{t} b_{\tau}\right) / 2\right) \lambda\left(x a_{\tau}\right)
$$

where $S[x]=\left(S_{v}\left[x_{v}\right]\right)_{v \in \mathbf{h}}$ and $\left(x_{1}, \ldots, x_{n}\right) a=\left(\sum_{i=1}^{n} x_{i} a_{i j}\right)_{j=1}^{n}$ for $\left(x_{1}, \ldots, x_{n}\right)$ $\in V_{\mathbf{h}}^{n}$ and $a \in\left(F_{\mathbf{h}}\right)_{n}^{n}$.

$$
\begin{equation*}
\left({ }^{\eta} \lambda\right)(x)=i^{r} \int_{Y} \lambda(y) \mathbf{e}_{\mathbf{h}}\left(-\sum_{i=1}^{n} S\left(x_{i}, y_{i}\right)\right) d^{S} y \tag{6}
\end{equation*}
$$

Here $Y=V_{\mathbf{h}}^{n}, d^{S} y$ is the Haar measure on $Y$ such that the measure of $\left(\sum_{i=1}^{q} \mathfrak{g}_{v} e_{i}\right)^{n}$ for each $v \in \mathbf{h}$ with a basis $\left\{e_{i}\right\}_{i=1}^{q}$ of $V$ over $F$ is $N\left(\mathcal{d}_{v}\right)^{-q n / 2}$ $\left|\operatorname{det}\left(S\left(e_{i}, e_{j}\right)\right)\right|_{v}^{n / 2}$, and $r=(n / 2) \sum_{v \in \mathbf{r}}\left(r_{v}-s_{v}\right)$ or $-n \sum_{v \in \mathbf{r}} s_{v}$ according as $q$ is even or odd.
(7) $J_{S}(-\eta, \eta z) J_{S}(\eta, z)=1$ and ${ }^{-\eta}\left({ }^{\eta} \lambda\right)=\lambda$.

Since $V$ has no fixed coordinate system, $\operatorname{det}(S)$ means the coset of $\operatorname{det}\left(S\left(e_{i}, e_{j}\right)\right)$ modulo $\left\{a^{2} \mid a \in F^{\times}\right\}$with $\left\{e_{i}\right\}_{i=1}^{n}$ as in (6) above. Clearly $\chi$ is well-defined. Thus our theorem is "coordinate-free" (as far as $V$ is concerned). In the proof, however, we use a matrix representation, and so hereafter we assume that $V=F^{q}$ and $S(x, y)={ }^{t} x S y$ for $x, y \in F_{n}^{q}$ with ${ }^{t} S=S \in G L_{q}(F)$. Then $A_{v} \in G L_{q}\left(F_{v}\right)$, $S_{v}={ }^{t} A_{v} I_{v} A_{v}$, and $T_{v}={ }^{t} A_{v} A_{v}$ for each $v \in \mathbf{a}$. In this setting, the formula of (6) can be written

$$
\begin{equation*}
\left({ }^{\eta} \lambda\right)(x)=i^{r} \mid N_{F / \mathbf{Q}}\left(\left.\operatorname{det}(S)\right|^{-n / 2} \int_{Y} \lambda(y) \mathbf{e}_{\mathbf{h}}\left(-\operatorname{tr}\left({ }^{t} x S y\right)\right) d y\right. \tag{A3.3}
\end{equation*}
$$

where $Y=\left(F_{n}^{q}\right)_{\mathbf{h}}$, and the measure of $\left(\mathfrak{g}_{v}\right)_{n}^{q}$ is $N\left(\mathfrak{p}_{v}\right)^{-n q / 2}$ for each $v \in \mathbf{h}$.
The proof of our theorem requires some preliminaries. We first define an embedding $\psi: \mathcal{H}^{(n)} \rightarrow \mathcal{H}^{(n q)}$ and an injective homomorphism $\alpha \mapsto \alpha_{S}$ of $G_{\mathbf{A}}^{(n)}$ into $G_{\mathbf{A}}^{(n q)}$ by $\psi(z)=\left(\psi_{v}\left(z_{v}\right)\right)_{v \in \mathbf{a}}$,

$$
\psi_{v}(x+i y)=S_{v} \otimes x+i T_{v} \otimes y, \quad \alpha_{S}=\left[\begin{array}{cc}
1_{q} \otimes a_{\alpha} & S \otimes b_{\alpha} \\
S^{-1} \otimes c_{\alpha} & 1_{q} \otimes d_{\alpha}
\end{array}\right]
$$

It can easily be verified that $\psi(\alpha(z))=\alpha_{S}(\psi(z))$ and

$$
\mu\left(\alpha_{S}, \psi(z)\right)_{v}=\left(A_{v} \otimes 1_{n}\right)^{-1} \operatorname{diag}\left[1_{r_{v}} \otimes \mu(\alpha, z)_{v}, 1_{s_{v}} \otimes \overline{\mu(\alpha, z)_{v}}\right]\left(A_{v} \otimes 1_{n}\right)
$$

for each $v \in \mathbf{a}$. From these we obtain immediately

$$
\begin{equation*}
j_{\alpha_{S}}(\psi(z))^{\mathbf{a}}=j_{\alpha}(z)^{q \mathbf{a}} \prod_{v \in \mathbf{a}} j_{\alpha}(z)_{v}^{-s_{v}}{\overline{j_{\alpha}(z)_{v}}}_{v}^{s_{v}} \tag{A3.4}
\end{equation*}
$$

To find a relationship between $J^{S}(\alpha, z)$ and $h\left(\alpha_{S}, \psi(z)\right)$, we take integral ideals $\mathfrak{b}$ and $\mathfrak{c}$ such that $\alpha_{S} \in D^{(n q)}\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]$ for every $\alpha \in D^{(n)}\left[2 \mathfrak{b} \mathfrak{d}^{-1}, 2 \mathfrak{c d}\right]$. Then $S \in G L_{q}\left(\mathfrak{g}_{v}\right)$ for every $v \nmid \mathfrak{b c}$.

A3.4. Lemma. For $\alpha \in G \cap P_{\mathbf{A}} D\left[2 \mathfrak{b d}^{-1}, 2 \mathfrak{c d}\right]$ we have

$$
h\left(\alpha_{S}, \psi(z)\right)=\chi_{\mathbf{a}}\left(\operatorname{det}\left(d_{\alpha}\right)\right) \chi^{*}\left(\operatorname{det}\left(d_{\alpha}\right) \operatorname{il}(\alpha)^{-1}\right) J^{S}(\alpha, z),
$$

where $\chi^{*}$ is the ideal character associated with $\chi$ of Theorem A3.3.
Proof. From Theorem A2.4 (1), (2), and (A3.4) we easily see that

$$
\begin{equation*}
h\left(\alpha_{S}, \psi(z)\right)=t J^{S}(\alpha, z) \quad \text { with } \quad t \in \mathbf{T} . \tag{A3.5}
\end{equation*}
$$

By Lemma A2.6 we have

$$
\lim _{z \rightarrow 0} h(\alpha, z) /|h(\alpha, z)|=\omega, \quad \lim _{z \rightarrow 0} h\left(\alpha_{S}, \psi(z)\right) /\left|h\left(\alpha_{S}, \psi(z)\right)\right|=\omega^{\prime}
$$

with $\omega=\omega\left(\delta^{-2} d_{\alpha}^{-1} c_{\alpha}\right)$ and $\omega^{\prime}=\omega\left(\delta^{-2} S^{-1} \otimes d_{\alpha}^{-1} c_{\alpha}\right)$. Let $\alpha \in P_{\mathbf{A}} \tau$ with $\tau \in$ $D\left[2 \mathfrak{b d} \boldsymbol{d}^{-1}, 2 \mathfrak{c d}\right]$; write simply $c$ and $d$ for $c_{\tau}$ and $d_{\tau}$. Clearly $d_{\alpha}^{-1} c_{\alpha}=d^{-1} c$. From (A3.2a) and (A3.4) we see that $\omega^{\prime}=t \chi_{\mathbf{a}}\left(\operatorname{det}\left(d_{\alpha}\right)\right) \xi$, where $\xi=1$ or $\xi=\omega^{q}$ according as $q$ is even or odd. If $v \nmid 2 \mathfrak{b c}$, then $S^{-1}={ }^{t} T_{v} \cdot \operatorname{diag}\left[s_{1}, \ldots, s_{q}\right] T_{v}$ with $T_{v} \in G L_{q}\left(\mathfrak{g}_{v}\right)$ and $s_{i} \in \mathfrak{g}_{v}^{\times}$. By Lemma A1.6 (1) we see that

$$
\gamma_{v}\left(\delta^{-2} S^{-1} \otimes d^{-1} c\right)=\prod_{i=1}^{q} \gamma_{v}\left(\delta^{-2} s_{i} d^{-1} c\right)=\gamma_{v}\left(\delta^{-2} d^{-1} c\right)^{q}\left(\frac{\operatorname{det}(S)}{\nu_{0}\left(\delta^{-1} d^{-1} c\right)_{v}}\right)
$$

By Lemma 1.11 (2) we have $\nu_{0}\left(\delta^{-1} d^{-1} c\right)=\nu_{0}\left(\delta^{-1} d_{\alpha}^{-1} c_{\alpha}\right)=\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{D}}(\alpha)^{-1}$. If $v \mid 2 \mathfrak{b c}$, then $|\operatorname{det}(d)|_{v}=1$, and hence both $\left(2^{-1} \delta^{-1} d^{-1} c\right)_{v}$ and $\left(2^{-1} \delta^{-1} S^{-1} \otimes d^{-1} c\right)_{v}$ are $v$-integral. Therefore (A1.5) shows that $\gamma_{v}=\gamma_{v}^{\prime}=1$. Notice that $\prod_{v \nmid 2} \omega_{v}^{q}=$ $\psi^{*}\left(\nu_{0}\left(\delta^{-1} d^{-1} c\right)\right)^{q / 2}$ by Lemma A1.6 (2), if $q$ is even, where $\psi^{*}$ is as in Proposition A2.7. Combining all these, we obtain our lemma.

Proof of Theorem A3.3. We are identifying $V$ with $F^{q}$ and $V^{n}$ with $F_{n}^{q}$, and hence $S_{v}[p]={ }^{t} p S_{v} p$ and $T_{v}[p]={ }^{t} p T_{v} p$ for $p \in\left(F_{n}^{q}\right)_{v}$. Define $\mathcal{A}: \mathcal{U}^{(n q)} \rightarrow \mathcal{U}^{(n q)}$ by $\mathcal{A}(u)_{v}=\left({ }^{t} A_{v} \otimes 1_{n}\right) u_{v}$ and also $\omega: F_{n q}^{1} \rightarrow F_{n}^{q}$ by $\omega\left(x_{1} \ldots x_{q}\right)={ }^{t}\left({ }^{t} x_{1} \ldots{ }^{t} x_{q}\right)$ for $x_{i} \in F_{n}^{1}$. Then a straightforward calculation shows that

$$
\begin{gather*}
g(u, z ; \lambda)=\theta^{(n q)}(\mathcal{A}(u), \psi(z) ; \lambda \circ \omega)  \tag{A3.6a}\\
\alpha_{S}(\mathcal{A}(u), \psi(z))=(\mathcal{A}(w), \psi(\alpha z)) \quad \text { if } \quad \alpha(u, z)=(w, \alpha z) \tag{A3.6b}
\end{gather*}
$$

For $\alpha \in G \cap \mathfrak{M}_{q}$ we can find an element $\xi \in \prod_{v \in \mathbf{a}} M p\left(\left(F_{n q}^{1}\right)_{v}\right)$ such that $\operatorname{pr}(\xi)=$ $\left(\alpha_{S}\right)_{\mathbf{a}}$ and $\prod_{v \in \mathbf{a}} g\left(\xi_{v}, \psi(z)_{v}\right)=J^{S}(\alpha, z)$. Then, for $\ell=\lambda \circ \omega$ Proposition A2.5 together with (A3.6a, b) shows that

$$
J^{S}(\alpha, z) g(u, z ; \lambda)=\theta^{(n q)}\left(\alpha_{S}(\mathcal{A}(u), \psi(z)) ; \ell^{\prime}\right)=g\left(\alpha(u, z) ; \lambda^{\prime}\right)
$$

with $\lambda^{\prime}=\ell^{\prime} \circ \omega^{-1}$. Putting $\lambda^{\prime}={ }^{\alpha} \lambda$, we obtain (0). Let $\mathfrak{k}$ be the conductor of $\chi$; let $E=D\left[2 \mathfrak{b} \mathfrak{d}^{-1}, 2 \mathfrak{c} \mathfrak{d}\right]$ and $E^{\prime}=\left\{\alpha \in E \mid \chi_{\mathfrak{k}}\left(\operatorname{det}\left(d_{\alpha}\right)\right)=1\right\}$, where $\chi_{\mathfrak{k}}=\prod_{v \mid \mathfrak{k}} \chi_{v}$. If $\alpha \in G \cap P_{\mathbf{A}} E$, then Lemma A3.4 and Proposition A2.5, combined with the above argument, show that

$$
\begin{equation*}
{ }^{\alpha} \lambda \circ \omega=\chi_{\mathbf{a}}\left(\operatorname{det}\left(d_{\alpha}\right)\right) \chi^{*}\left(\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{o}}(\alpha)^{-1}\right) \cdot{ }^{\beta}(\lambda \circ \omega) \tag{*}
\end{equation*}
$$

with $\beta=\alpha_{S}$. In particular, ${ }^{\alpha} \lambda \circ \omega={ }^{\beta}(\lambda \circ \omega)$ if $\alpha \in G \cap E^{\prime}$. This together with Theorem A2.4, (6) shows that ${ }^{\alpha} \lambda=\lambda$ for every $\alpha$ in a congruence subgroup $\Gamma$ of $G$ depending on $\lambda$. Take an open subgroup $D$ of $E^{\prime}$ so that $\Gamma \supset G \cap D$. We take $D \subset C^{\theta}$ if $q$ is odd. By strong approximation, given $\sigma \in \mathfrak{M}_{q}$, we can take $\alpha \in G$ so that $\operatorname{pr}(\sigma) \in \alpha D$. Then $\alpha \in \mathfrak{M}_{q}$. Define ${ }^{\sigma} \lambda$ to be ${ }^{\alpha} \lambda$. It is then easy to verify that this is well-defined and has property (2) for even $q$, and also properties (3) and (4). To prove (5), let $\operatorname{pr}(\sigma)=\tau \in P_{\mathbf{h}}$. Given $\lambda \in \mathcal{S}\left(\left(F_{n}^{q}\right)_{\mathbf{h}}\right)$, take $\alpha \in G$ so that $\tau \in \alpha E^{\prime}$ and that ${ }^{\sigma} \lambda={ }^{\alpha} \lambda$. Then we easily see that

$$
\begin{equation*}
\chi_{\mathbf{a}}\left(\operatorname{det}\left(d_{\alpha}\right)\right) \chi^{*}\left(\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{v}}(\alpha)^{-1}\right)=\chi\left(\operatorname{det}\left(a_{\tau}\right)\right) \tag{**}
\end{equation*}
$$

Let $\ell=\lambda \circ \omega$ and $\beta=\alpha_{S}$. From $\left(^{*}\right)$ and $\left({ }^{* *}\right)$ we obtain ${ }^{\beta} \ell=\chi\left(\operatorname{det}\left(a_{\tau}\right)\right)\left({ }^{\alpha} \lambda\right) \circ \omega$. Changing $E^{\prime}$ for a smaller group if necessary, we may assume that $\tau_{S}$ and $\alpha_{S}$ have the same effect on $\ell$. Then we obtain
$\left.{ }^{* * *}\right) \quad \quad \quad\left(\operatorname{det}\left(a_{\tau}\right)\right)\left({ }^{\alpha} \lambda\right) \circ \omega={ }^{\varphi} \ell \quad$ with $\operatorname{pr}(\varphi)=\tau_{S}$.
Taking $\varphi$ to be $r_{P}\left(\tau_{S}\right)$, from Theorem A2.4 (8) and (A2.3a) we obtain

$$
\varphi_{\ell} \ell(y)=\left|\operatorname{det}\left(1_{q} \otimes a_{\tau}\right)\right|_{\mathbf{A}}^{1 / 2} \mathbf{e}_{\mathbf{h}}\left(y\left(S \otimes a_{\tau} \cdot{ }^{t} b_{\tau}\right) \cdot{ }^{t} y / 2\right) \ell\left(y\left(1_{q} \otimes a_{\tau}\right)\right),
$$

which combined with $\left(^{* * *}\right)$ proves (5). To prove (6), or rather (A3.3), we first observe that (A3.5) is valid for $\alpha=\eta$. Now the first formula of Lemma A2.6 shows that both $h(\eta, z)$ and $h\left(\eta_{S}, \psi(z)\right)$ are positive if $z \in \mathbf{R i}$ with $\mathbf{i}$ of that lemma, and hence $i^{r} \cdot h\left(\eta_{S}, \psi(z)\right)=J^{S}(\eta, z)$ with $r$ as in (6). Then from (A3.6a, b) we obtain ${ }^{\eta} \lambda \circ \omega=i^{r} .{ }^{\gamma}(\lambda \circ \omega)$ with $\gamma=\eta_{S}$. Observe that $\gamma$ belongs to the set $P_{\mathbf{A}} C^{*}$ of Theorem A2.4 (7) (of degree nq). Therefore, by (A2.3c) we have

$$
{ }^{\gamma} \ell(x)=\left|\operatorname{det}\left(c_{\gamma}\right)_{\mathbf{h}}\right|_{\mathbf{A}}^{1 / 2} \int_{Y} \ell\left(y c_{\gamma}\right) \mathbf{e}_{\mathbf{h}}\left(x b_{\gamma} \cdot{ }^{t} c_{\gamma} \cdot{ }^{t} y\right) d y \quad\left(Y=\left(F_{n q}^{1}\right)_{\mathbf{h}}\right)
$$

which can easily be transformed to (A3.3). To prove (2) when $q$ is odd, first let $\sigma \in \mathfrak{M}$ and $\operatorname{pr}(\rho) \in C^{\theta}$. With an open normal subgroup $D$ of $C^{\theta}$, take $\alpha \in$ $G \cap \operatorname{pr}(\sigma) D$ and $\beta \in G \cap \operatorname{pr}(\rho) D$. Then $\alpha \beta \in G \cap \operatorname{pr}(\sigma \rho) D$. Take $D$ so small that ${ }^{\sigma}\left({ }^{\beta} \lambda\right)={ }^{\alpha}\left({ }^{\beta} \lambda\right),{ }^{\beta} \lambda={ }^{\rho} \lambda$, and ${ }^{(\sigma \rho)} \lambda={ }^{(\alpha \beta)} \lambda$. Since $\beta \in G \cap C^{\theta}$, (A3.2b) shows that $J^{S}(\alpha \beta, z)=J^{S}(\alpha, \beta z) J^{S}(\beta, z)$, and hence we obtain ${ }^{(\alpha \beta)} \lambda={ }^{\alpha}\left({ }^{\beta} \lambda\right)$ from (0). Thus ${ }^{(\sigma \rho)} \lambda={ }^{(\alpha \beta)} \lambda={ }^{\alpha}\left({ }^{\beta} \lambda\right)={ }^{\sigma}\left({ }^{\beta} \lambda\right)={ }^{\sigma}\left({ }^{\rho} \lambda\right)$. Next let $\operatorname{pr}(\tau) \in P_{\mathbf{A}}$ and $\sigma \in \mathfrak{M}$. Then $\sigma=\pi \rho$ with $\operatorname{pr}(\pi) \in P_{\mathbf{A}}$ and $\operatorname{pr}(\rho) \in C^{\theta}$. From ( ${ }^{* * *}$ ) we see that ${ }^{(\tau \pi)} \zeta={ }^{\tau}\left({ }^{\pi} \zeta\right)$ for every $\zeta \in \mathcal{S}\left(\left(F_{m}^{q}\right)_{\mathbf{h}}\right)$. Therefore ${ }^{(\tau \sigma)} \lambda={ }^{(\tau \pi \rho)} \lambda={ }^{(\tau \pi)}\left({ }^{\rho} \lambda\right)=$ ${ }^{\tau}\left({ }^{\pi}\left({ }^{\rho} \lambda\right)\right)={ }^{\tau}\left({ }^{(\pi \rho)} \lambda\right)={ }^{\tau}\left({ }^{\sigma} \lambda\right)$. This proves (2). To prove (1), it is sufficient to show that ${ }^{\alpha} \lambda$ for $\alpha \in G \cap \mathfrak{M}_{q}$ is independent of $\left\{A_{v}\right\}$. If $q$ is even, this follows from (5) and (6), since $P$ and $\eta$ generate $G$. Suppose $q$ is odd. By (2) and (5), it is sufficient to show that ${ }^{\alpha} \lambda$ for $\alpha \in G \cap C^{\theta}$ is independent of $\left\{A_{v}\right\}$. Now $\alpha^{m} \in D\left[2 \mathfrak{b d}^{-1}, 2 \mathfrak{c d}\right]$ for some positive integer $m$. Therefore Lemma A3.4 shows
that $h\left(\alpha_{S}, \psi(z)\right)=t J^{S}(\alpha, z)$ with a root of unity $t$. Then ${ }^{\alpha} \lambda \circ \omega=t^{-1} \cdot \gamma(\lambda \circ \omega)$ with $\gamma=\alpha_{S}$. Now the set of all $\left\{A_{v}\right\}$ is not necessarily connected, but it can easily be shown that the set of all $\left\{T_{v}\right\}$ is connected. Since $h\left(\alpha_{S}, \psi(z)\right)$ is continuous in $\left\{T_{v}\right\}$, we see that $t$ does not depend on $\left\{A_{v}\right\}$. This proves (1). Finally, to prove (7), we recall that $h(-\eta, z)=h(\eta, z)$ as stated in Theorem A2.4 (4). Therefore from (A2.24) we easily see that $h(-\eta, z) h(\eta, z)=1$. This combined with (A3.2a) proves the first half of (7). Then the second half of (7) follows from (0).

A3.5. We are going to introduce a series involving harmonic polynomials on $V^{n}$. Given a finite-dimensional complex vector space $W$ and $0 \leq a \in \mathbf{Z}$, we denote by $\mathfrak{S}_{a}(W)$ the vector space of all $\mathbf{C}$-valued homogeneous polynomial functions on $W$ of degree $a$. We then denote by $\mathcal{P}_{a}\left(\mathbf{C}_{m}^{q}\right)$ the vector subspace of $\mathfrak{S}_{a}\left(\mathbf{C}_{m}^{q}\right)$ spanned by the functions $p$ satisfying the condition

$$
\begin{equation*}
\sum_{i=1}^{q} \partial^{2} p / \partial x_{i h} \partial x_{i k}=0 \quad \text { for every } h \text { and } k \tag{A3.7a}
\end{equation*}
$$

where $x=\left(x_{i h}\right)$ is a variable on $\mathbf{C}_{m}^{q}$. For instance, we can take $p(x)=\varphi\left({ }^{t} \rho x\right)$ with $\varphi \in \mathfrak{S}_{a}\left(\mathbf{C}_{m}^{m}\right)$ and $\rho \in \mathbf{C}_{m}^{q}$ satisfying the condition

$$
\begin{equation*}
{ }^{t} \rho_{h} \rho_{k}=0 \text { whenever } \partial^{2} \varphi / \partial y_{h i} \partial y_{k j} \neq 0 \text { for some } i \text { and } j, \tag{A3.7b}
\end{equation*}
$$

where $\rho_{h}$ denotes the $h$-th column of $\rho$, and $y=\left(y_{h i}\right)$ is a variable on $\mathbf{C}_{m}^{m}$.
A3.6. Lemma. Let $\omega(x)=\exp \left(\sum_{h, k=1}^{m} \sum_{i=1}^{q} c_{h k} x_{i h} x_{i k}\right)$ for $x \in \mathbf{R}_{m}^{q}$ with $c_{h k}=$ $c_{k h} \in \mathbf{C}$. Then $[p(D)(\omega \psi)](0)=[p(D) \psi](0)$ for every $p \in \mathcal{P}_{a}\left(\mathbf{C}_{m}^{q}\right)$ and every $C^{\infty}$ function $\psi$ in $x$, where $D$ is the $(q \times m)$-matrix whose $(i, h)$-entry is $\partial / \partial x_{i h}$.

Proof. We first observe that if $\alpha$ is a polynomial in $n$ variables $y_{1}, \ldots, y_{n}$ and $\alpha_{i}=\partial \alpha / \partial y_{i}$, then

$$
\begin{equation*}
\left[\alpha\left(\partial / \partial y_{1}, \ldots, \partial / \partial y_{n}\right)\left(y_{i} \beta\right)\right](0)=\left[\alpha_{i}\left(\partial / \partial y_{1}, \ldots, \partial / \partial y_{n}\right) \beta\right](0) \tag{*}
\end{equation*}
$$

for every $i$ and every $C^{\infty}$ function $\beta$. This is completely elementary. Now our lemma is trivial if $a=0$. Assume that it is true for degree $<a$ and that $a>0$. We have $a p(x)=\sum_{i, h} x_{i h} p_{i h}(x)$ with $p_{i h}=\partial p / \partial x_{i h}$, and hence

$$
\begin{aligned}
{[a p(D)(\omega \psi)](0) } & =\sum_{i, h}\left[p_{i h}(D) \partial / \partial x_{i h}(\omega \psi)\right](0) \\
& =\sum_{i, h}\left[p_{i h}(D)\left(\omega \cdot \partial \psi / \partial x_{i h}\right)\right](0)+\sum_{i, h, k} 2 c_{h k}\left[p_{i h}(D)\left(x_{i k} \omega \psi\right)\right](0)
\end{aligned}
$$

Since $p_{i h} \in \mathcal{P}_{a-1}\left(\mathbf{C}_{m}^{q}\right)$, by our induction assumption the first sum on the last line is $\sum_{i, h}\left[p_{i h}(D) \partial \psi / \partial x_{i h}\right](0)$, which is $[a p(D) \psi](0)$. By (*) the second sum equals $\sum_{i, h, k} 2 c_{h k}\left[\left(\partial p_{i h} / \partial x_{i k}\right)(D)(\omega \psi)\right](0)$, which is 0 by (A3.7a). This completes the proof.

A3.7. Coming back to the space $V$ and the form $S$, for each $v \in$ a put

$$
\begin{align*}
& X_{v}^{+}=\left\{x \in V_{v} \mid\left(A_{v} x\right)_{i}=0 \quad \text { for } \quad i>r_{v}\right\}  \tag{A3.8a}\\
& X_{v}^{-}=\left\{x \in V_{v} \mid\left(A_{v} x\right)_{i}=0 \quad \text { for } \quad i \leq r_{v}\right\} \tag{A3.8b}
\end{align*}
$$

where $y_{i}$ for $y \in \mathbf{R}^{q}$ means the $i$-th component of $y$. Clearly $V_{v}=X_{v}^{+} \oplus X_{v}^{-}$. For $y \in V_{v}$, we denote by $y^{+}$and $y^{-}$the projections of $y$ to $X_{v}^{+}$and $X_{v}^{-}$.

Given $m$ and $m^{\prime}$ in $\mathbf{Z}^{\mathbf{a}}$ whose components are all nonnegative, we denote by $\mathcal{P}_{m, m^{\prime}}\left(V^{n}\right)$ the vector space over $\mathbf{C}$ spanned by all functions $p$ on $V_{\mathbf{a}}^{n}=\prod_{v \in \mathbf{a}} V_{v}^{n}$ of the form

$$
\begin{equation*}
p(x)=\prod_{v \in \mathbf{a}} p_{v}\left(A_{v} x_{v}^{+}\right) p_{v}^{\prime}\left(A_{v} x_{v}^{-}\right) \tag{A3.9}
\end{equation*}
$$

with $p_{v} \in \mathcal{P}_{m_{v}}\left(\mathbf{C}_{n}^{r_{v}}\right), p_{v}^{\prime} \in \mathcal{P}_{m_{v}^{\prime}}\left(\mathbf{C}_{n}^{s_{v}}\right)$, where $A_{v} x_{v}^{ \pm}=\left(A_{v} x_{v 1}^{ \pm}, \ldots, A_{v} x_{v n}^{ \pm}\right)$for $x_{v}=\left(x_{v 1}, \ldots, x_{v n}\right)$ with $x_{v i} \in V_{v}$. Write $p$ of (A3.9) as $p=\left(p_{v}, p_{v}^{\prime}\right)$. We let every element $\mu$ of $G L_{n}(\mathbf{C})^{\mathbf{a}}$ act $\mathbf{C}$-linearly on $\mathcal{P}_{m, m^{\prime}}\left(V^{n}\right)$ by defining $\mu p=$ $\left(\mu_{v} p_{v}, \bar{\mu}_{v} p_{v}^{\prime}\right)$ with $(\nu s)(y)=s(y \nu)$ for $s=p_{v}, p_{v}^{\prime}$, and $\nu=\mu_{v}$ or $\bar{\mu}_{v}$. Notice that this action is compatible with both (A3.7a) and (A3.7b).

Now for $z \in \mathcal{H}, \lambda \in \mathcal{S}\left(V_{\mathrm{h}}^{n}\right)$, and $p \in \mathcal{P}_{m, m^{\prime}}\left(V^{n}\right)$ we consider a series

$$
\begin{equation*}
f(z ; \lambda, p)=\sum_{\xi \in V^{n}} \lambda\left(\xi_{\mathbf{h}}\right) p\left(\xi_{\mathbf{a}}\right) \Phi(\xi ; 0, z) . \tag{A3.10}
\end{equation*}
$$

A3.8. Theorem. The notation ${ }^{\alpha} \lambda$ being as in Theorem A3.3, we have

$$
J^{S}(\alpha, z)^{-1} f\left(\alpha(z) ;{ }^{\alpha} \lambda,{ }^{t} \mu(\alpha, z)^{-1} p\right)=f(z ; \lambda, p)
$$

for every $\alpha \in G \cap \mathfrak{M}_{q}$.
Proof. Write the variable $u$ in the form ${ }^{t} u_{v}=\left(u_{v 1}^{1}, \ldots, u_{v 1}^{n}, \ldots, u_{v q}^{1}, \ldots, u_{v q}^{n}\right)$. We can then define a differential operator $B=\prod_{v \in \mathbf{a}} p_{v}\left(D_{v}\right) p_{v}^{\prime}\left(D_{v}^{\prime}\right)$ on $\mathcal{U}$, where $D_{v}=\left(\partial / \partial u_{v i}^{j}\right)$ with $1 \leq i \leq r_{v}, 1 \leq j \leq n, D_{v}^{\prime}=\left(\partial / \partial u_{v i}^{j}\right)$ with $r_{v}<i \leq q, 1 \leq$ $j \leq n, E_{v}=\left(\partial / \partial a_{v i}^{j}\right), E_{v}^{\prime}=\left(\partial / \partial b_{v i}^{j}\right)$ with $1 \leq i \leq q, 1 \leq j \leq n$. Employing Lemma A3.6 we can easily verify that $[B g(u, z ; \lambda)]_{u=0}=(2 \pi i)^{N} f(z ; \lambda, p)$, where $N=\sum_{v \in \mathbf{a}}\left(m_{v}+m_{v}^{\prime}\right)$. Therefore we obtain our assertion by applying $B$ to the equality of Theorem A3.3 (0).

We can associate with the above $f$ a function $f_{\mathbf{A}}(x ; \lambda, p)$ with a variable $x$ on $G_{\mathbf{A}}$ or $M_{\mathbf{A}}$, according as $q$ is even or odd, by

$$
\begin{equation*}
f_{\mathbf{A}}(\alpha w ; \lambda, p)=J^{S}(w, \mathbf{i})^{-1} f\left(w(\mathbf{i}) ;{ }^{w} \lambda,{ }^{t} \mu(w, \mathbf{i})^{-1} p\right) \tag{A3.11}
\end{equation*}
$$

for $\alpha \in G$ and $w \in \mathfrak{M}_{q}$, where $\mathbf{i}$ is as in Lemma A2.6; we take $\operatorname{pr}(w) \in C^{\theta}$ if $q$ is odd. This is well-defined by virtue of Theorem A3.8. Now we have

A3.9. Proposition. For every $\alpha \in G$ and $y \in \mathfrak{M}_{q}$ such that $y(\mathbf{i})=\mathbf{i}$ and that $\operatorname{pr}(y) \in C^{\theta}$ if $q$ is odd, we have

$$
f_{\mathbf{A}}(\alpha x y ; \lambda, p)=J^{S}(y, \mathbf{i})^{-1} f_{\mathbf{A}}\left(x ;{ }^{y} \lambda,{ }^{t} \mu(y, \mathbf{i})^{-1} p\right)
$$

Proof. Given $x$, take $\beta \in G$ and $w \in \mathfrak{M}_{q}$ so that $x=\beta w$ and $\operatorname{pr}(w) \in C^{\theta}$. Then

$$
\begin{aligned}
f_{\mathbf{A}}(\alpha x y, \lambda, p) & =f_{\mathbf{A}}(w y, \lambda, p) \\
& =J^{S}(w y, \mathbf{i})^{-1} f\left(w(\mathbf{i}),{ }^{(w y)} \lambda,{ }^{t} \mu(w y, \mathbf{i})^{-1} p\right) \\
& =J^{S}(y, \mathbf{i})^{-1} J^{S}(w, \mathbf{i})^{-1} f\left(w(\mathbf{i}),{ }^{\left.w\left({ }^{y} \lambda\right),{ }^{t} \mu(w, \mathbf{i})^{-1} \cdot{ }^{t} \mu(y, \mathbf{i})^{-1} p\right)}\right. \\
& =J^{S}(y, \mathbf{i})^{-1} f_{\mathbf{A}}\left(x,{ }^{y} \lambda,{ }^{t} \mu(y, \mathbf{i})^{-1} p\right) .
\end{aligned}
$$

A3.10. Proposition. The notation being as in Proposition A2.10, we have $J^{S}\left(\sigma^{*},-z^{\rho}\right)=\overline{J^{S}(\sigma, z)}$ and ${ }^{\sigma^{*}}\left(\lambda^{*}\right)=\left({ }^{\sigma} \lambda\right)^{*}$ for every $\sigma \in \mathfrak{M}_{q}$ and $\lambda \in \mathcal{S}\left(V_{\mathbf{h}}^{n}\right)$, where $\lambda^{*}$ is defined by $\lambda^{*}(x)=\overline{\lambda(-x)}$.

Proof. The first equality follows from Proposition A2.10 and (A3.2a) immediately. By virtue of strong approximation and (4) of Theorem A3.3, it is sufficient to prove the second assertion when $\sigma \in G \cap \mathfrak{M}_{q}$, in which case the desired fact follows from (0) of Theorem A3.3, since $g\left(\bar{u},-z^{\rho} ; \lambda^{*}\right)=\overline{g(u, z ; \lambda)}$.

A3.11. Remark. (I) Define an algebraic group $O(S)$ by

$$
O(S)=\{\alpha \in G L(V) \mid S(\alpha x, \alpha x)=S(x, x)\}
$$

Fixing $\left(A_{v}\right)_{v \in \mathbf{a}}$ as above, put $A_{v}^{\alpha}=A_{v} \alpha_{v}$ for every $\alpha \in O(S)_{\mathbf{a}}$. Then we can define our series $g$ and $f$ with $A_{v}^{\alpha}$ in place of $A_{v}$. Thus $g$ and $f$ are essentially parametrized by $O(S)_{\mathbf{a}}$.
(II) If $r \leq \operatorname{Min}(q, m)$, we easily see that the subdeterminants of $x \in \mathbf{C}_{m}^{q}$ of degree $r$ define elements of $\mathcal{P}_{r}\left(\mathbf{C}_{m}^{q}\right)$. In particular, if $r_{v}=q=n$, we can take $p_{v}\left(A_{v} x_{v}^{+}\right)=\operatorname{det}\left(A_{v} x_{v}\right)$ in (A3.9). In this case $\mu_{v} p_{v}=\operatorname{det}\left(\mu_{v}\right) p_{v}$.
(III) Take $q \geq n$ and $S=T \in G L_{q}(F)$ with $r_{v}=q$ and $s_{v}=0$ for every $v \in \mathbf{a}$; take also a subset $\mathbf{a}^{\prime}$ of $\mathbf{a}$. For $x \in \mathbf{C}_{n}^{q}$ and $v \in \mathbf{a}^{\prime}$ let $p_{v}(x)$ be the determinant of the first $n$ rows of $x$; let $p_{v}=1$ if $v \notin \mathbf{a}^{\prime}$. We have clearly

$$
f(z ; \lambda, p)=\sum_{\xi \in F_{n}^{q}} \lambda\left(\xi_{\mathbf{h}}\right) p\left(\xi_{\mathbf{a}}\right) \mathbf{e}_{\mathbf{a}}^{n}\left(2^{-1} \cdot{ }^{t} \xi S \xi z\right)
$$

Then Theorem A3.8 shows that this is an element of $\mathcal{M}_{k}$ with $k=(q / 2) \mathbf{a}+\mathbf{a}^{\prime}$. It is a cusp form if $\mathbf{a}^{\prime} \neq \varnothing$. Indeed, by Lemma $7.5, G$ is generated by $G \cap \mathfrak{M}_{q}$, and hence Theorem A3.8 shows that the tranform of $f(z ; \lambda, p)$ by an element of $G$ (understood in the sense of (10.12) if $k \notin \mathbf{Z}^{\mathbf{a}}$ ) is of the form $f\left(z ; \lambda^{\prime}, p\right.$ ) with some $\lambda^{\prime}$. If $\mathbf{a}^{\prime} \neq \varnothing$, we have clearly $f\left(z ; \lambda^{\prime}, p\right)=\sum_{h \in S} c(h) \mathbf{e}_{\mathbf{a}}^{n}(h z)$ with $c(h) \neq 0$ only for $\operatorname{det}(h) \neq 0$; thus $f(z ; \lambda, p)$ is a cusp form. To find $\lambda$ such that $f(z ; \lambda, p) \neq 0$, take $\xi \in \mathfrak{g}_{n}^{q}$ so that $p\left(\xi_{\mathbf{a}}\right) \neq 0$ and put $X=\left\{\left.\xi^{\prime} \in \mathfrak{g}_{n}^{q}\right|^{t} \xi^{\prime} S \xi^{\prime}={ }^{t} \xi S \xi\right\}$. Then $X$ is a finite set. Therefore we can easily find $\lambda \in \mathcal{S}\left(\left(F_{n}^{q}\right)_{\mathbf{h}}\right)$ such that $\lambda(\xi) \neq 0$ and $\lambda\left(\xi^{\prime}\right)=0$ if $\xi \neq \xi^{\prime} \in X$ or $\xi^{\prime} \notin \mathfrak{g}_{n}^{q}$. Then clearly $f(z ; \lambda, p) \neq 0$.

Let us now derive from Theorem A3.3 explicit formulas for ${ }^{\alpha} \lambda$ in two forms convenient in applications. For simplicity, we put ${ }^{\alpha} \lambda={ }^{\beta} \lambda$ if $\alpha=\operatorname{pr}(\beta)$ with $\beta \in \mathfrak{M}_{q}, q$ odd. This is meaningful in view of Theorem A3.3 (3). Now our first formula is:

A3.12. Lemma. Let $\mathfrak{k}$ be the conductor of $\chi$ of Theorem A3.3, and $\chi^{*}$ the ideal character associated with $\chi$. For $\lambda \in \mathcal{S}\left(V_{\mathbf{h}}^{n}\right)$ put

$$
U_{\lambda}=\left\{\left.\sigma \in D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{k d}\right]\right|^{\sigma} \lambda=\lambda \text { and } \operatorname{det}\left(d_{\sigma}\right)_{v} \equiv 1\left(\bmod \mathfrak{k}_{v}\right) \text { for every } v \mid \mathfrak{k}\right\}
$$

Then, for every $\alpha \in \operatorname{diag}\left[p,{ }^{t} p^{-1}\right] U_{\lambda}$ with $p \in G L_{n}\left(F_{\mathbf{h}}\right)$ we have

$$
\left({ }^{\alpha} \lambda\right)(x)=\left|\operatorname{det}(p)_{\mathbf{h}}\right|_{\mathbf{A}}^{q / 2} \chi_{\mathbf{h}}(\operatorname{det}(p)) \lambda(x p)
$$

where $x p$ is as in Theorem A3.3 (5). In particular, if such an $\alpha$ belongs to $G$, then

$$
\left({ }^{\alpha} \lambda\right)(x)=\chi_{\mathbf{a}}\left(\operatorname{det}\left(d_{\alpha}\right)\right) \chi^{*}\left(\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{d}}(\alpha)^{-1}\right) N\left(\mathrm{il}_{\mathfrak{D}}(\alpha)\right)^{q / 2} \lambda(x p)
$$

Proof. Let $\alpha=\tau \sigma$ with $\tau=\operatorname{diag}\left[p,{ }^{t} p^{-1}\right]$ and $\sigma \in U_{\lambda}$. By Theorem A3.3 $(2,5)$ we have $\left({ }^{\alpha} \lambda\right)(x)=\left({ }^{\tau} \lambda\right)(x)=\left|\operatorname{det}(p)_{\mathbf{h}}\right|_{\mathbf{A}}^{q / 2} \chi_{\mathbf{h}}(\operatorname{det}(p)) \lambda(x p)$. Suppose $\alpha \in G$. Since ${ }^{t} p d_{\alpha}=d_{\sigma}$, we have

$$
\chi_{\mathbf{a}}\left(\operatorname{det}\left(d_{\alpha}\right)\right) \chi^{*}\left(\operatorname{det}\left(d_{\alpha}\right) \mathrm{il}_{\mathfrak{d}}(\alpha)^{-1}\right)=\chi_{\mathbf{h}}(\operatorname{det}(p)) \chi_{\mathfrak{k}}\left(\operatorname{det}\left(d_{\sigma}^{-1}\right)\right)=\chi_{\mathbf{h}}(\operatorname{det}(p))
$$

which completes the proof.

A3.13. Proposition. Given $\lambda \in \mathcal{S}\left(V_{\mathrm{h}}^{n}\right)$, let $M$ be a g-lattice in $V^{n}$ such that $\lambda(x+u)=\lambda(x)$ for every $u \in M$. Further let $\mathfrak{x}, \mathfrak{y}$, and $\mathfrak{z}$ be fractional ideals of $F$ with the following properties:
(i) $2\left(1+\delta_{i j}\right)^{-1} S\left(x_{i}, x_{j}\right) \in \mathfrak{x}$ for every $i, j$ and every $x \in V^{n}$ such that $\lambda(x) \neq 0$.
(ii) $2\left(1+\delta_{i j}\right)^{-1} S\left(y_{i}, y_{j}\right) \in \mathfrak{y}$ for every $i, j$ and every $y \in M^{\prime}$, where

$$
M^{\prime}=\left\{y \in V^{n} \mid \sum_{i=1}^{n} S\left(x_{i}, y_{i}\right) \in \mathfrak{d}^{-1} \text { for every } \quad x \in M\right\}
$$

(iii) $\lambda(x a)=\lambda(x)$ for every $a \in \prod_{v \in \mathrm{~h}} G L_{n}\left(\mathfrak{g}_{v}\right)$ such that $a_{v}-1 \in\left(\mathfrak{z}_{v}\right)_{n}^{n}$ for every $v \in \mathbf{h}$, where $x a$ is as in Theorem A3.3 (5).
Then with $\chi$ and $\mathfrak{k}$ as in Theorem A3.3 and Lemma A3.12 we have

$$
\left({ }^{\gamma} \lambda\right)(x)=\chi_{\mathfrak{k}}\left(\operatorname{det}\left(a_{\gamma}\right)\right) \lambda\left(x\left(a_{\gamma}\right)_{\mathfrak{z}}\right) \quad \text { for every } \gamma \in B
$$

where $\left(a_{\gamma}\right)_{\mathfrak{z}}$ is the projection of $a_{\gamma}$ to $\prod_{v \mid \mathfrak{z}} G L_{n}\left(F_{v}\right), B=D\left[2 \mathfrak{d}^{-1} \mathfrak{y}^{-1}, 2 \mathfrak{d}^{-1} \mathfrak{y}^{-1} \cap\right.$ $\left.\left\{2^{-1} \mathfrak{d x}(\mathfrak{k} \cap \mathfrak{z})\right\}\right]$ if $q$ is even, and $B=D\left[2 \mathfrak{d}^{-1} \mathfrak{a}, 2^{-1} \mathfrak{d} \mathfrak{a}^{-1} \mathfrak{b}\right], \mathfrak{a}=\mathfrak{x}^{-1} \cap \mathfrak{g}, \mathfrak{b}=$ $\mathfrak{k} \cap \mathfrak{z} \cap 4 \mathfrak{a} \cap 4 \mathfrak{d}^{-2} \mathfrak{a} \mathfrak{y}^{-1}$ if $q$ is odd.

We first prove:
A3.14. Lemma. Let $\mathfrak{b}$ and $\mathfrak{c}$ be fractional ideals in $F$ such that $\mathfrak{b c}$ is integral, and $\mathfrak{a}$ an integral ideal such that $\mathfrak{a} \subset \mathfrak{b} \cap \mathfrak{c}$. Further let $E(\mathfrak{b})=G_{\mathbf{a}} \prod_{v \in \mathfrak{h}} E_{v}(\mathfrak{b})$ and $E^{\prime}(\mathfrak{c})=G_{\mathbf{a}} \prod_{v \in \mathbf{h}} E_{v}^{\prime}(\mathfrak{c})$, where $E_{v}(\mathfrak{b})$ resp. $E_{v}^{\prime}(\mathfrak{c})$ denotes the set of all elements of $G_{v}$ of the form $\left[\begin{array}{cc}1_{n} & b \\ 0 & 1_{n}\end{array}\right]$ resp. $\left[\begin{array}{cc}1_{n} & 0 \\ c & 1_{n}\end{array}\right]$ with $b \in\left(\mathfrak{b}_{v}\right)_{n}^{n}$ resp. $c \in\left(\mathfrak{c}_{v}\right)_{n}^{n}$. Then $D[\mathfrak{b}, \mathfrak{c}]$ is generated by $E(\mathfrak{b}), E^{\prime}(\mathfrak{c})$, and $D[\mathfrak{a}, \mathfrak{a}]$.

Proof. Since $D_{v}[\mathfrak{a}, \mathfrak{a}]=D_{v}[\mathfrak{b}, \mathfrak{c}]$ if $v \nmid \mathfrak{a}$, it is sufficient to show that $D_{v}[\mathfrak{b}, \mathfrak{c}]$ is generated by $E_{v}(\mathfrak{b}), E_{v}^{\prime}(\mathfrak{c})$, and $D_{v}[\mathfrak{a}, \mathfrak{a}]$. If $v \nmid \mathfrak{b} \mathfrak{c}$, then $D_{v}[\mathfrak{b}, \mathfrak{c}]$ is conjugate to $D_{v}[\mathfrak{g}, \mathfrak{g}]=S p\left(n, \mathfrak{g}_{v}\right)$, and hence our assertion follows from a well-known fact that $S p\left(n, \mathfrak{g}_{v}\right)$ is generated by $E_{v}(\mathfrak{g}), E_{v}^{\prime}(\mathfrak{g})$, and $\operatorname{diag}\left[a,{ }^{t} a^{-1}\right]$ with $a \in G L_{n}\left(\mathfrak{g}_{v}\right)$. If $v \mid \mathfrak{b c}$ and $\left[\begin{array}{ll}a & b \\ c & d\end{array}\right] \in D_{v}[\mathfrak{b}, \mathfrak{c}]$, then ${ }^{t} a d-1 \prec(\mathfrak{b} \mathfrak{c})_{v}$, and hence $a \in G L_{n}\left(\mathfrak{g}_{v}\right)$, and

$$
\left[\begin{array}{ll}
a & b \\
c & d
\end{array}\right]=\left[\begin{array}{cc}
1 & 0 \\
c a^{-1} & 1
\end{array}\right]\left[\begin{array}{cc}
a & 0 \\
0 & { }^{t} a^{-1}
\end{array}\right]\left[\begin{array}{cc}
1 & a^{-1} b \\
0 & 1
\end{array}\right]
$$

which proves our lemma.
Proof of Proposition A3.13. We use the matrix representation as in the proof of Theorem A3.3. Let $\alpha \in P_{\mathbf{A}}$ with $a_{\alpha}=1_{n}$. Then Theorem A3.3 (5) shows that ${ }^{\alpha} \lambda(x)=\lambda(x) \mathbf{e}_{\mathbf{h}}\left(\operatorname{tr}\left({ }^{t} x S x b_{\alpha}\right) / 2\right)$. Therefore ${ }^{\alpha} \lambda=\lambda$ if $\alpha \in E\left(2 \mathfrak{d}^{-1} \mathfrak{x}^{-1}\right)$. Put $\beta=\eta^{-1} \alpha \eta$ and $\lambda^{\prime}={ }^{\eta} \lambda$. Substituting $y+z$ for $y$ in (A3.3), we find that $\lambda^{\prime}(x)=\mathbf{e}_{\mathbf{h}}\left(\operatorname{tr}\left({ }^{t} x S z\right)\right) \lambda^{\prime}(x)$ for every $z \in M$, and hence $\lambda^{\prime}(x) \neq 0$ only if $x \in M^{\prime}$. By (ii) this means that $\lambda^{\prime}(x) \neq 0$ only if ${ }^{t} x S x$ has entries in $\mathfrak{y}$. Therefore ${ }^{\alpha} \lambda^{\prime}=\lambda^{\prime}$ if $\alpha \in E\left(2 \mathfrak{d}^{-1} \mathfrak{y}^{-1}\right)$. Suppose $\beta \in E^{\prime}\left(2 \mathfrak{d} \cap 2 \mathfrak{d}^{-1} \mathfrak{y}^{-1}\right)$. Since $\beta \in C^{\theta}$ and $\alpha \in P_{\mathbf{A}}$, we have ${ }^{\eta}\left({ }^{\beta} \lambda\right)={ }^{(\eta \beta)} \lambda={ }^{(\alpha \eta)} \lambda={ }^{\alpha}\left({ }^{\eta} \lambda\right)={ }^{\eta} \lambda$, and hence ${ }^{\beta} \lambda=\lambda$. By Lemma A3.12 the expected formula for ${ }^{\gamma} \lambda$ is true for $\gamma \in D[\mathfrak{e}, \mathfrak{e}]$ with a suitable ideal $\mathfrak{e}$. We have seen that it is also true for $\gamma \in E\left(2 \mathfrak{d}^{-1} \mathfrak{x}^{-1}\right) \cup E^{\prime}\left(2 \mathfrak{d} \cap 2 \mathfrak{d}^{-1} \mathfrak{y}^{-1}\right)$, which together with Lemma A3.14 proves our proposition for odd $q$, since ${ }^{\delta}\left({ }^{\varepsilon} \lambda\right)={ }^{\delta \varepsilon} \lambda$ at least for $\delta, \varepsilon \in C^{\theta}$. If $q$ is even, the associativity is true for all $\delta, \varepsilon \in G_{\mathbf{A}}$, and so ${ }^{\beta} \lambda=\lambda$ for $\beta \in E\left(2 \mathfrak{d}^{-1} \mathfrak{y}^{-1}\right)$. Therefore we can take $B$ in the form stated in our proposition.

The following lemma, though unnecessary in the present book, is of independent interest, and so we give here a proof.

A3.15. Lemma. The notation being as in Lemma A3.14, let $T(\mathfrak{b})=G \cap E(\mathfrak{b})$ and $T^{\prime}(\mathfrak{c})=G \cap E^{\prime}(\mathfrak{c})$. Then $\Gamma[\mathfrak{b}, \mathfrak{c}]$ is generated by $T(\mathfrak{b}), T^{\prime}(\mathfrak{c})$, and $\Gamma[\mathfrak{a}, \mathfrak{a}]$.

Proof. Let $X$ be an open normal subgroup of $D[\mathfrak{b}, \mathfrak{c}]$ contained in $D[\mathfrak{a}, \mathfrak{a}]$. Given $\alpha \in \Gamma[\mathfrak{b}, \mathfrak{c}]$, Lemma A3.14 allows us to take $u_{1}, \ldots, u_{m}$ in $E(\mathfrak{b}) \cup E^{\prime}(\mathfrak{c}) \cup$ $D[\mathfrak{a}, \mathfrak{a}]$ so that $\alpha=u_{1} \cdots u_{m}$. By strong approximation, $u_{i} \in \beta_{i} X$ with some $\beta_{i} \in G$. If $u_{i} \in D[\mathfrak{a}, \mathfrak{a}]$, then $\beta_{i} \in G \cap D[\mathfrak{a}, \mathfrak{a}]=\Gamma[\mathfrak{a}, \mathfrak{a}]$. If $u_{i} \in E(\mathfrak{b})$, we can take $\beta_{i}$ from $T(\beta)$, and similarly if $u_{i} \in E^{\prime}(\mathfrak{c})$, we can take $\beta_{i}$ from $T^{\prime}(\mathfrak{c})$. Then $\alpha=u_{1} \cdots u_{m} \in \beta_{1} \cdots \beta_{m} X$, and hence $\alpha=\beta_{1} \cdots \beta_{m} \gamma$ with $\gamma \in G \cap X \subset \Gamma[\mathfrak{a}, \mathfrak{a}]$, which completes the proof.

A3.16. We now consider the special class of theta series by taking $n=q$. Thus we put $W=F_{n}^{n}$, and identify it with $V^{n}$. For $z \in \mathfrak{H}^{\mathbf{a}}, \lambda \in \mathcal{S}\left(W_{\mathbf{h}}\right)$, a totally positive symmetric element $\tau$ of $W$, and $\mu \in \mathbf{Z}^{\mathbf{a}}$ such that $0 \leq \mu_{v} \leq 1$ for every $v \in \mathbf{a}$, put

$$
\begin{equation*}
\theta(z, \lambda)=\sum_{\xi \in W} \lambda\left(\xi_{\mathbf{h}}\right) \operatorname{det}(\xi)^{\mu} \mathbf{e}_{\mathbf{a}}^{n}\left({ }^{t} \xi \tau \xi z\right) \tag{A3.12}
\end{equation*}
$$

This is a special case of the function defined by (A3.10). Indeed, let $S(x, y)=$ $2 \cdot{ }^{t} x \tau y$ for $x, y \in V=F_{1}^{n}$; we can take $p(\xi)=\operatorname{det}(\xi)^{\mu}$ for $\xi \in W$ as explained in Remark A3.11 (II). Then $\theta$ of (A3.12) can be obtained as $f(z ; \lambda, p)$ of (A3.10). We now put $\mathfrak{M}_{n}=\mathfrak{M}$ with $\mathfrak{M}$ of (A2.17) if $n$ is odd and $\mathfrak{M}_{n}=G_{\mathbf{A}}$ if $n$ is even. Putting $l=\mu+(n / 2)$ a, we define a factor of automorphy $J(\alpha, z)$ for $\alpha \in \mathfrak{M}_{n}$ by

$$
J(\alpha, z)= \begin{cases}j_{\alpha}^{l}(z) & \text { if } n \text { is even }  \tag{A3.13}\\ h(\alpha, z)^{n} j_{\alpha}^{\mu}(z) & \text { if } n \text { is odd }\end{cases}
$$

From Theorem A3.8 we obtain

$$
\begin{equation*}
\theta\left(\alpha z,^{\alpha} \lambda\right)=J(\alpha, z) \theta(z, \lambda) \quad \text { for every } \quad \alpha \in G \cap \mathfrak{M}_{n} \tag{A3.14}
\end{equation*}
$$

Moreover, for each $\lambda$, our function $\theta(z, \lambda)$ is an element of $\mathcal{M}_{l}$. (See $\S 6.10$ for the definition of $\mathcal{M}_{l}$ if $l \notin \mathbf{Z}^{\mathbf{a}}$.) Now, by the principle of (A3.11) we can associate with the above $\theta$ a function $\theta_{\mathbf{A}}^{\prime}(x, \lambda)$ with a variable $x$ on $G_{\mathbf{A}}$ or $M_{\mathbf{A}}$, according as $n$ is even or odd, by

$$
\begin{equation*}
\theta_{\mathbf{A}}^{\prime}(x, \lambda)=J(w, \mathbf{i})^{-1} \theta\left(w(\mathbf{i}),{ }^{w} \lambda\right) \tag{A3.15}
\end{equation*}
$$

for $x=\alpha w$ with $\alpha \in G$ and $w \in \mathfrak{M}_{n}$; we take $\operatorname{pr}(w) \in C^{\theta}$ if $n$ is odd. This is well-defined. It should be noted that $\theta_{\mathbf{A}}^{\prime}$ is the function associated with $\theta$ by the principle of $\S 16.6$ only if $J=j^{l}$, which is not necessarily true if $n$ is odd. In the rest of this section we put

$$
S=\left\{\left.x \in W\right|^{t} x=x\right\}
$$

We no longer use $S$ of $\S 3.1$ in the rest of Section A3, but reinstate it in Section A4.
As a special case of what we said in Remark A3.11 (III) we have
(A3.16) $\theta(z, \lambda)$ is a cusp form if $\mu \neq 0$.
A3.17. Proposition. Let $\chi$ be the Hecke character of $F$ corresponding to the extension $F\left(\operatorname{det}(2 \tau)^{1 / 2}\right) / F$ or $F\left((-1)^{n / 4} \operatorname{det}(\tau)^{1 / 2}\right) / F$ according as $n$ is odd or even. Then

$$
\begin{align*}
\theta_{\mathbf{A}}^{\prime}\left(r_{P}\left[\begin{array}{cc}
q & s \hat{q} \\
0 & \widehat{q}
\end{array}\right], \lambda\right) & =\chi(\operatorname{det}(q)) \operatorname{det}(q)_{\mathbf{a}}^{\mu}|\operatorname{det}(q)|_{\mathbf{A}}^{n / 2}  \tag{A3.17}\\
& \cdot \sum_{\xi \in W} \lambda\left(\xi_{\mathbf{h}} q\right) \operatorname{det}(\xi)^{\mu} \mathbf{e}_{\mathbf{a}}^{n}\left(\mathbf{i} \cdot{ }^{t} q \cdot{ }^{t} \xi \tau \xi q\right) \mathbf{e}_{\mathbf{A}}^{n}\left({ }^{t} \xi \tau \xi s\right)
\end{align*}
$$

for every $q \in G L_{n}(F)_{\mathbf{A}}$ and $s \in S_{\mathbf{A}}$, where $r_{p}$ is the identity map of $G_{\mathbf{A}}$ onto itself if $n$ is even. Moreover if $\beta=r_{P}(\operatorname{diag}[r, \widehat{r}]) w$ with $\beta \in G, r \in G L_{n}(F)_{\mathbf{h}}$, and $w \in \mathfrak{M}_{n}, \operatorname{pr}(w) \in D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]$, then

$$
\begin{align*}
J\left(\beta, \beta^{-1} z\right) \theta\left(\beta^{-1} z, \lambda\right)= & \chi_{\mathbf{h}}(\operatorname{det}(r))|\operatorname{det}(r)|_{\mathbf{A}}^{n / 2}  \tag{A3.18}\\
& \cdot \sum_{\xi \in W}\left({ }^{w} \lambda\right)\left(\xi_{\mathbf{h}} r\right) \operatorname{det}(\xi)^{\mu} \mathbf{e}_{\mathbf{a}}^{n}\left({ }^{t} \xi \tau \xi z\right) .
\end{align*}
$$

Proof. Given $x \in \mathfrak{M}_{n}$, take $\alpha \in G$ and $w$ as in (A3.15). Then $\alpha \in \mathfrak{M}_{n}$. Put $z=w(\mathbf{i})$. By (A3.14) and (A3.15) we have $\theta_{\mathbf{A}}^{\prime}(x, \lambda)=J(w, \mathbf{i})^{-1} \theta\left(z,{ }^{w} \lambda\right)=$ $J(w, \mathbf{i})^{-1} J(\alpha, z)^{-1} \theta\left(\alpha z,{ }^{\alpha}\left({ }^{w} \lambda\right)\right)$. By (A3.2b) and Theorem A3.3 (2) we have $J(\alpha w, \mathbf{i})=J(\alpha, z) J(w, \mathbf{i})$ and ${ }^{\alpha}\left({ }^{w} \lambda\right)={ }^{x} \lambda$ since $\operatorname{pr}(w) \in C^{\theta}$. Therefore

$$
\begin{equation*}
\theta_{\mathbf{A}}^{\prime}(x, \lambda)=J(x, \mathbf{i})^{-1} \theta\left(x(\mathbf{i}),{ }^{x} \lambda\right) \quad \text { if } \quad x \in \mathfrak{M}_{n} \tag{A3.19}
\end{equation*}
$$

Take $x=r_{P}\left[\begin{array}{cc}q & s \widehat{q} \\ 0 & \widehat{q}\end{array}\right]$. By Theorem A3.3 (5) we have

$$
\left({ }^{x} \lambda\right)(y)=\left|\operatorname{det}(q)_{\mathbf{h}}\right|_{\mathbf{A}}^{n / 2} \chi_{\mathbf{h}}(\operatorname{det}(q)) \mathbf{e}_{\mathbf{h}}^{n}\left({ }^{t} y \tau y s\right) \lambda(y q) \quad\left(y \in W_{\mathbf{h}}\right)
$$

This combined with (A3.12) and (A3.19) gives (A3.17). To prove (A3.18), take the element $x$ so that $\operatorname{pr}(x) \in P_{\mathbf{a}}$ and $x(\mathbf{i})=z$. Since $\beta^{-1} z=w^{-1} x(\mathbf{i})$ with $\beta$ as in (A3.18) and $G_{\mathbf{a}}$ acts trivially on $\mathcal{S}\left(W_{\mathbf{h}}\right)$, (A3.15) shows that $J\left(\beta, \beta^{-1} z\right) \theta\left(\beta^{-1} z, \lambda\right)$ $=J\left(\beta, \beta^{-1} z\right) J\left(w^{-1} x, \mathbf{i}\right) \theta_{\mathbf{A}}^{\prime}\left(w^{-1} x,{ }^{w} \lambda\right)=J\left(\beta w^{-1} x, \mathbf{i}\right) \theta_{\mathbf{A}}^{\prime}\left(\beta w^{-1} x,{ }^{w} \lambda\right)=J(g, \mathbf{i})$ $\cdot \theta_{\mathbf{A}}^{\prime}\left(g,{ }^{w} \lambda\right)$, where $g=r_{P}\left[\begin{array}{cc}r q & r s \widehat{q} \\ 0 & \widehat{r} \widehat{q}\end{array}\right]$. Thus we obtain (A3.18) from (A3.17).

A3.18. Fixing an integral ideal $\mathfrak{e}$ in $F$, we put

$$
\begin{gathered}
R=\prod_{v \in \mathbf{h}}\left(\mathfrak{g}_{v}\right)_{n}^{n}\left(\subset W_{\mathbf{h}}\right), \quad E_{v}=G L_{n}\left(\mathfrak{g}_{v}\right), \quad E_{v}^{\prime}=\left\{x \in E_{v} \mid x-1 \prec \mathfrak{e}_{v}\right\} \\
R^{\prime}=\left\{x \in R \mid x_{v} \in E_{v}^{\prime} \quad \text { for every } v \mid \mathfrak{e}\right\}, \quad R^{*}=R^{\prime} \cdot W_{\mathbf{a}}\left(\subset W_{\mathbf{A}}\right)
\end{gathered}
$$

We take $\mu$ and $\tau$ as above and a Hecke character $\omega$ of $F$ such that

$$
\begin{equation*}
\omega_{\mathbf{a}}(-1)^{n}=(-1)^{n\|\mu\|}, \quad\|\mu\|=\sum_{v \in \mathbf{a}} \mu_{v}, \quad \text { if } 2 \in \mathfrak{e} \tag{A3.20}
\end{equation*}
$$

and denote the conductor of $\omega$ by $\mathfrak{f}$. Taking an element $p$ of $G L_{n}(F)_{\mathbf{h}}$, we define a series $\theta$ by

$$
\begin{equation*}
\theta(z)=\sum_{\xi \in W \cap p R^{*}} \omega_{\mathbf{a}}(\operatorname{det}(\xi)) \omega^{*}\left(\operatorname{det}\left(p^{-1} \xi\right) \mathfrak{g}\right) \operatorname{det}(\xi)^{\mu} \mathbf{e}_{\mathbf{a}}^{n}\left({ }^{t} \xi \tau \xi z\right) . \tag{A3.21}
\end{equation*}
$$

Here it is understood that $\omega_{\mathbf{a}}(b) \omega^{*}(b \mathfrak{a})$ for $b=0$ and a fractional ideal $\mathfrak{a}$ denotes $\omega^{*}(\mathfrak{a})$ or 0 according as $\mathfrak{f}=\mathfrak{g}$ or $\mathfrak{f} \neq \mathfrak{g}$. For $\alpha \in \mathfrak{M}_{n}$ we define $j_{\alpha}^{l}(z)$ by (16.17).

The ideal $\mathfrak{e}$ is needed only for some technical reasons; the series of (A3.21) is most natural when $\mathfrak{e}=\mathfrak{g}$. Notice that $\theta$ is identically equal to 0 if condition (A3.20) is not satisfied, which can happen only if $n$ is odd.

A3.19. Proposition. Let $\rho_{\tau}$ be the Hecke character of $F$ corresponding to the extension $F\left(c^{1 / 2}\right) / F$ with $c=(-1)^{[n / 2]} \operatorname{det}(2 \tau) ;$ put $\mathfrak{f}^{\prime}=\mathfrak{f} \cap \mathfrak{e}$ and $\omega^{\prime}=\omega \rho_{\tau}$. Then there exist a fractional ideal $\mathfrak{b}$ and an integral ideal $\mathfrak{c}$ such that $\mathfrak{c} \subset \mathfrak{e}$, the conductor of $\omega^{\prime}$ divides $\mathfrak{c}, D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right] \subset D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]$ if $n$ is odd, and

$$
\begin{equation*}
\theta(\gamma z)=\omega_{\mathfrak{c}}^{\prime}\left(\operatorname{det}\left(a_{\gamma}\right)\right) j_{\gamma}^{l}(z) \theta(z) \quad \text { for every } \quad \gamma \in G \cap C \tag{A3.22}
\end{equation*}
$$

where $C=\left\{x \in D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right] \mid a_{x}-1 \prec \mathfrak{e}\right\}$. Moreover, if $\beta \in G \cap \operatorname{diag}[r, \widehat{r}] C$ with $r \in G L_{n}(F)_{\mathbf{h}}$, then

$$
\begin{align*}
j_{\beta}^{l}\left(\beta^{-1} z\right) \theta\left(\beta^{-1} z\right)=\omega^{\prime}(\operatorname{det}(r))^{-1} \omega_{\mathrm{c}}^{\prime}\left(\operatorname{det}\left(d_{\beta} r\right)\right)|\operatorname{det}(r)|_{\mathbf{A}}^{n / 2}  \tag{A3.23}\\
\sum_{\xi \in W \cap p R^{*} r^{-1}} \omega_{\mathbf{a}}(\operatorname{det}(\xi)) \omega^{*}\left(\operatorname{det}\left(\xi p^{-1} r\right) \mathfrak{g}\right) \operatorname{det}(\xi)^{\mu} \mathbf{e}_{\mathbf{a}}^{n}\left({ }^{t} \xi \tau \xi z\right)
\end{align*}
$$

In particular, suppose that ${ }^{t} g \cdot 2 \tau g \in \mathfrak{x}$ for every $g \in p L_{0}$ and ${ }^{t} h(2 \tau)^{-1} h \in 4 t^{-1}$ for every $h \in \widehat{p} L_{0}$ with fractional ideals $\mathfrak{x}$ and $\mathfrak{t}$, where $L_{0}=\mathfrak{g}_{1}^{n}$; let $\mathfrak{h}$ be the conductor of $\rho_{\tau}$. Then we can take $(\mathfrak{b}, \mathfrak{c})=\left(2^{-1} \mathfrak{d} \mathfrak{x}, \mathfrak{h} \cap \mathfrak{f}^{\prime} \cap \mathfrak{x}^{-1} \mathfrak{f}^{\prime 2} \mathfrak{t}\right)$ if $n$ is even and $(\mathfrak{b}, \mathfrak{c})=\left(2^{-1} \mathfrak{d} \mathfrak{a}^{-1}, \mathfrak{h} \cap \mathfrak{f}^{\prime} \cap 4 \mathfrak{a} \cap \mathfrak{a} \mathfrak{f}^{\prime 2} \mathfrak{t}\right)$ if $n$ is odd, where $\mathfrak{a}=\mathfrak{x}^{-1} \cap \mathfrak{g}$.

Proof. For $\beta$ and $r$ as above, we have $\operatorname{il}_{\mathfrak{D}}(\beta)=\operatorname{det}(r)^{-1} \mathfrak{g}$; thus by Proposition A2.7,

$$
\begin{equation*}
h(\beta, z)^{2}=\psi_{2}\left(\operatorname{det}\left(d_{\beta} r\right)\right) \psi_{\mathbf{h}}(\operatorname{det}(r)) j_{\beta}(z)^{\mathbf{a}} \tag{A3.24}
\end{equation*}
$$

where $\psi$ is as in that proposition and $\psi_{2}=\prod_{v \mid 2} \psi_{v}$. It follows that

$$
\begin{equation*}
J(\beta, z)=\left\{\psi_{2}\left(\operatorname{det}\left(d_{\beta} r\right)\right) \psi_{\mathbf{h}}(\operatorname{det}(r))\right\}^{[n / 2]} j_{\beta}^{l}(z) \quad \text { if } n \notin 2 \mathbf{Z} \tag{A3.25}
\end{equation*}
$$

In particular $J(\alpha, z)=\psi_{2}\left(\operatorname{det}\left(d_{\alpha}\right)\right)^{[n / 2]} j_{\alpha}^{l}(z)$ if $\alpha \in G \cap D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]$ and $n \notin 2 \mathbf{Z}$. Now define $\lambda \in \mathcal{S}\left(W_{\mathbf{h}}\right)$ by $\lambda(x)=\omega_{\mathbf{h}}\left(\operatorname{det}(p)^{-1}\right) \lambda^{\prime}\left(p^{-1} x\right), \lambda^{\prime}(x)=\prod_{v \in \mathbf{h}} \lambda_{v}^{\prime}\left(x_{v}\right)$, with

$$
\lambda_{v}^{\prime}(y)= \begin{cases}1 & \text { if } y \in\left(\mathfrak{g}_{v}\right)_{n}^{n}, v \nmid \mathfrak{f}^{\prime} \\ \omega_{v}\left(\operatorname{det}(y)^{-1}\right) & \text { if } y \in E_{v}^{\prime}, v \mid \mathfrak{f}^{\prime} \\ 0 & \text { otherwise }\end{cases}
$$

Then the function $\theta(z, \lambda)$ of (A3.12) coincides with $\theta$ of (A3.21). Notice that $\lambda(x a)=\omega_{f^{\prime}}(\operatorname{det}(a))^{-1} \lambda(x)$ for every $a \in \prod_{v \in \mathbf{h}} E_{v}^{\prime}$. Applying Proposition A3.13 to the present $\lambda$, we find a fractional ideal $\mathfrak{b}$ and an integral ideal $\mathfrak{c}$ such that $D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right] \subset D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]$ if $n$ is odd, the conductors of $\chi$ and $\omega$ divide $\mathfrak{c}$, and that

$$
\begin{equation*}
\left({ }^{w} \lambda\right)(x)=\chi_{\mathfrak{k}}\left(\operatorname{det}\left(a_{w}\right)\right) \lambda\left(x\left(a_{w}\right)_{\mathfrak{f}^{\prime}}\right) \quad \text { for every } \quad w \in D\left[\mathfrak{b}^{-1}, \mathfrak{b} \mathbf{c}\right], \tag{A3.26}
\end{equation*}
$$

where $\chi$ is as in Proposition A3.17, and $\mathfrak{k}$ is its conductor. This combined with (A3.14), (A3.18), and (A3.25) proves our assertions up to formula (A3.23). Let $\mathfrak{h}^{\prime}$ be the conductor of $\chi$. Then $\mathfrak{h}^{\prime}=\mathfrak{h}$ if $n$ is even and $\mathfrak{h} \cap 4 \mathfrak{g}=\mathfrak{h}^{\prime} \cap 4 \mathfrak{g}$ if $n$ is odd. Therefore Proposition A3.13 gives our assertion on $(\mathfrak{b}, \mathfrak{c})$ as stated above. (In fact, the ideal $\mathfrak{y}$ of that proposition is $4\left(\mathfrak{d}^{2} \mathfrak{f}^{\prime 2} \mathfrak{t}\right)^{-1}$ in the present case.)

A3.20. Example. Take $\theta$ of (A3.21) with trivial $\omega, \mu=\mathbf{a}, \mathfrak{e}=\mathfrak{g}, \tau=2^{-1} 1_{n}$, $p=1_{n}$, and $n \in 4 \mathbf{Z}$. Then $l=(1+(n / 2))$ a and $\rho_{\tau}$ is trivial; thus we can take $\mathfrak{x}=\mathfrak{g}$ and $\mathfrak{t}=4 \mathfrak{g}$ in Proposition A3.19, so that $\mathfrak{b}=2^{-1} \mathfrak{d}$ and $\mathfrak{c}=4 \mathfrak{g}$. Therefore $\theta \in \mathcal{M}_{l}(\Gamma)$ with $\Gamma=G \cap D\left[2 \mathfrak{d}^{-1}, 2 \mathfrak{d}\right]$. By (A3.16), $\theta$ is a cusp form. If ${ }^{t} \xi \xi={ }^{t} \xi_{1} \xi_{1}$, then $\operatorname{det}(\xi)= \pm \operatorname{det}\left(\xi_{1}\right)$, and so $\operatorname{det}(\xi)^{\mathbf{a}}=\operatorname{det}\left(\xi_{1}\right)^{\mathbf{a}}$ if $[F: \mathbf{Q}]$ is even. If $[F: \mathbf{Q}]$ is odd, then ${ }^{t}(\beta \xi)(\beta \xi)={ }^{t} \xi \xi$ and $\operatorname{det}(\beta \xi)^{\mathbf{a}}=-\operatorname{det}(\xi)^{\mathbf{a}}$ with $\beta=\operatorname{diag}\left[-1,1_{n-1}\right]$. Thus $\theta$ is nonzero if and only if $[F: \mathbf{Q}]$ is even.

## A4. The constant term of a theta series at each cusp depends only on the genus

A4.1. Our purpose is to prove the fact stated in the title, and to discuss its consequences in connection with Siegel's theory of quadratic forms. Our setting is the same as in $\S A 3.1$. For simplicity we use matrix expressions; thus taking ${ }^{t} S=S \in G L_{q}(F)$, we put $W=F_{n}^{q}$,

$$
O(S)=\left\{\left.\alpha \in G L_{q}(F)\right|^{t} \alpha S \alpha=S\right\}
$$

and $S[x]={ }^{t} x S x$ for $x \in W$; we also identify $W$ with $V^{n}$. For $\lambda \in \mathcal{S}\left(W_{\mathrm{h}}\right)$ we view $\lambda$ also as a function on $W_{\mathbf{A}}$ by putting $\lambda(x)=\lambda\left(x_{\mathbf{h}}\right)$.

A4.2. Theorem. Let $\alpha \in G$ and $\lambda \in \mathcal{S}\left(W_{\mathbf{h}}\right)$. Define $g(u, z ; \lambda)$ by (A3.0a) with $W$ in place of $V^{n}$. Let $t(z)=\zeta \prod_{v \in \mathbf{a}} j_{\alpha}(z)_{v}^{(q / 2)-s_{v}}\left|j_{\alpha}(z)_{v}\right|^{s_{v}}$ with $\zeta \in \mathbf{T}$ and any choice of the branch of $j_{\alpha}(z)_{v}^{(q / 2)-s_{v}}$. (If $\alpha \in \mathfrak{M}_{q}$, we can take $J^{S}(\alpha, z)$ to be $t(z)$.) Then the following assertions hold:
(1) There exists $\lambda_{1} \in \mathcal{S}\left(W_{\mathbf{h}}\right)$ determined by the equation

$$
\begin{equation*}
t(z)^{-1} g(\alpha(u, z) ; \lambda)=g\left(u, z ; \lambda_{1}\right) \tag{*}
\end{equation*}
$$

(2) Define $\mu \in \mathcal{S}\left(W_{\mathbf{h}}\right)$ by $\mu(x)=\lambda(\gamma x)$ with an element $\gamma \in O(S)_{\mathbf{h}}$; define $\mu_{1}$ by taking $\mu$ in place of $\lambda$ in ( ${ }^{*}$ ). Then $\mu_{1}(x)=\lambda_{1}(\gamma x)$.

Proof. Given $\lambda$ and $\alpha \in G \cap \mathfrak{M}_{q}$, take $\ell$ so that ${ }^{\alpha} \ell=\lambda$ and put $\ell=\lambda_{1}$. Then from Theorem A3.3 (0) we obtain

$$
J^{S}(\alpha, z)^{-1} g(\alpha(u, z) ; \lambda)=g\left(u, z ; \lambda_{1}\right)
$$

Now $G \cap \mathfrak{M}_{q}$ contains $P$ and $\eta$. By Lemma 7.5 every element of $G$ is a product of finitely many elements in $P \cup\{\eta\}$, since $\eta^{-1}=-\eta$ and $-1 \in P$. Now if the assertion of (1) is true for $(\alpha, t)$ and ( $\alpha^{\prime}, t^{\prime}$ ), then we easily see that it is true for $\left(\alpha \alpha^{\prime}, t^{\prime \prime}\right)$, where $t^{\prime \prime}(z)=\zeta t\left(\alpha^{\prime} z\right) t^{\prime}(z)$ with any $\zeta \in \mathbf{T}$. Therefore we obtain (1). As for (2), for the same reason it is sufficient to prove it when $\alpha \in P$ or $\alpha=-\eta$ with $t(z)=J^{S}(\alpha, z)$. If $\alpha \in P$, taking $\beta=\alpha^{-1}$, we have $\left(^{*}\right)$ with $\lambda_{1}={ }^{\beta} \lambda$. By Theorem A3.3 (5) we have

$$
\left({ }^{\beta} \lambda\right)(x)=\varepsilon(\beta) \mathbf{e}_{\mathbf{h}}\left(\operatorname{tr}\left(S[x] a_{\beta} \cdot{ }^{t} b_{\beta}\right) / 2\right) \lambda\left(x a_{\beta}\right)
$$

with a constant $\varepsilon(\beta)$ that depends only on $\beta$ and $S$. Then clearly $\left({ }^{\beta} \lambda\right)(\gamma x)=$ $\left({ }^{\beta} \mu\right)(x)$, which is the desired relation for $\alpha \in P$. Next take $\alpha=-\eta$. Then ( ${ }^{*}$ ) holds with $\lambda_{1}={ }^{\eta} \lambda$ by Theorem A3.3 (7). Now ${ }^{\eta} \lambda$ is given by (A3.3). For $\gamma$ and $\mu$ as above, we have $d(\gamma y)=d y$, and hence we can easily derive from (A3.3) that $\left({ }^{\eta} \mu\right)(x)=\left({ }^{\eta} \lambda\right)(\gamma x)$ as desired. This completes the proof.

A4.3. Theorem. (1) Suppose that $n=1$; given $\lambda \in \mathcal{S}\left(W_{\mathbf{h}}\right)$ and $\gamma \in O(S)_{\mathbf{h}}$, define $\mu \in \mathcal{S}\left(W_{\mathbf{h}}\right)$ by $\mu(x)=\lambda(\gamma x)$. Then $g(0, z ; \lambda)-g(0, z ; \mu)$ is rapidly decreasing at the cusps of $G$.
(2) Suppose moreover that $S$ is totally definite. Put

$$
\begin{equation*}
f(z, \lambda)=\sum_{g \in W} \lambda(g) \mathbf{e}_{\mathbf{a}}(S[g] z / 2) \quad\left(z \in \mathfrak{H}_{1}^{\mathbf{a}}\right) \tag{A4.1}
\end{equation*}
$$

If $\mu$ is as above, then $f(z, \lambda)-f(z, \mu)$ is a cusp form.
Proof. Put $\nu=\lambda-\mu$; then clearly $g(0, z ; \lambda)-g(0, z ; \mu)=g(0, z ; \nu)$, and if $\alpha$ and $t(z)$ are as in Theorem A4.2, then

$$
\begin{equation*}
t(z)^{-1} g(0, \alpha z ; \nu)=g\left(0, z ; \nu_{1}\right) \tag{A4.2}
\end{equation*}
$$

with $\nu_{1}=\lambda_{1}-\mu_{1}$. By Theorem A4.2 (2), we have $\mu_{1}(x)=\lambda_{1}(\gamma x)$, so that $\nu_{1}(0)=0$. Since $n=1$, the definition of the series $g$ shows that the function of (A4.2) must be rapidly decreasing. This proves (1). Assertion (2) is merely a special case of (1).

A4.4. Theorem. Let $L$ be a $\mathfrak{g}$-lattice in $F^{q}$, and let

$$
f(z, L)=\sum_{g \in L} \mathbf{e}_{\mathbf{a}}(S[g] z / 2) \quad\left(z \in \mathfrak{H}_{1}^{\mathbf{a}}\right)
$$

Then the following assertions hold:
(1) $f(z, L)-f(z, \gamma L)$ is a cusp form for every $\gamma \in O(S)_{\mathbf{h}}$.
(2) Let $\left\{L_{i}\right\}_{i=1}^{h}$ be a complete set of representatives for the classes of lattices in the genus of $L$ with respect to $O(S)$. Put $e_{i}=\#\left\{\alpha \in O(S) \mid \alpha L_{i}=L_{i}\right\}$ and

$$
\begin{equation*}
p(z)=\left(\sum_{i=1}^{h} e_{i}^{-1}\right)^{-1} \sum_{i=1}^{h} e_{i}^{-1} f\left(z, L_{i}\right) \tag{A4.3}
\end{equation*}
$$

Then $f(z, L)-p(z)$ is a cusp form.
Proof. Clearly $h(z, L)=h(z, \lambda)$ if we take $\lambda(x)=\prod_{v \in \mathbf{h}} \lambda_{v}\left(x_{v}\right)$ with the characteristic function of $L_{v}$ as $\lambda_{v}$. Therefore (1) follows from Theorem A4.3 (2). Assertion (2) follows immediately from (1).

Now $p$ of (A4.3) equals $F(S, z)$ with the function $F$ defined by Siegel in $[\mathrm{Si}$, I, p.372, (78); p.542, (129)]. He proved (2) of Theorems A4.4 when $F=\mathbf{Q}$ in [ $\mathrm{Si}, \mathrm{I}, \mathrm{p} .376$ ]. The number $\sum_{i=1}^{h} e_{i}^{-1}$ is the mass of the genus of $L$ in his sense. If we put $f(z, L)=\sum_{b \in F} r(b, L) \mathbf{e}_{\mathbf{a}}(b z / 2)$ and $p(z)=\sum_{b \in F} r_{0}(b) \mathbf{e}_{\mathbf{a}}(b z / 2)$, then $r(b, L)=\#\{g \in L \mid S[g]=b\}$ and $r_{0}(b)=\left(\sum_{i=1}^{h} e_{i}^{-1}\right)^{-1} \sum_{i=1}^{h} e_{i}^{-1} r\left(b, L_{i}\right)$. Thus $r_{0}(b)$ is the weighted average of the numbers of representations of $b$ by $S$, for which Siegel gave his product formula.

We can also show that $p$ is an Eisenstein series by means of the Siegel-Weil formula. Siegel gave this fact in [Si, I, p.373, Satz 3; p.543, Satz III] for $n>4$. The case of an arbitrary $n$ is explained in [S99, Section 5].

In Sections A2 through A4 we have assumed that the basic filed $F$ is totally real, but we can actually treat the case of an arbitrary number field. Indeed, such a theory was presented in [S93], which gives at least the generalizations of the results up to Lemma A3.15 without assuming $F$ to be totally real. Also, in [S97, Section A7] we treated theta series of a hermitian form in a similar fashion. In the next section we will give some more results complementary to this theory. Generalizations or analogues of Theorems A4.2, A4.3, and A4.4 can be proved for theta series over an arbitrary number field and also in the hermitian case by the same methods.

## A5. Theta series of a hermitian form

This section concerns Case UT. Thus $K$ is a CM-field; see $\S 3.5$. For $y \in K_{\mathbf{A}}^{\times}$and $\mu \in \mathbf{Z}^{\mathbf{a}}$ we note that $|y|^{\mu}=\prod_{v \in \mathbf{a}}\left|y_{v}\right|^{\mu_{v}}$, where $\left|y_{v}\right|$ is the standard absolute value in $\mathbf{C}$, not its square. For $b \in K^{\times}$, for example, we have $\left|b_{\mathbf{h}}\right|_{K}^{-1}=\left|b_{\mathbf{a}}\right|_{K}=|b|^{2 \mathbf{a}}$.

A5.1. Lemma. Let $\varphi_{0}$ be a Hecke character of $F$ such that $\varphi_{0}(x)=x^{\nu}|x|^{-\nu}$ for $x \in F_{\mathbf{a}}^{\times}$with $\nu \in \mathbf{Z}^{\mathbf{a}}$. Then there exists a Hecke character $\varphi$ of $K$ such that $\varphi=\varphi_{0}$ on $F_{\mathbf{A}}^{\times}, \varphi(y)=y^{\nu}|y|^{-\nu}$ for $y \in K_{\mathbf{a}}^{\times}$, and the conductor of $\varphi$ divides a power of the conductor of $\varphi_{0}$.

Proof. Let $a$ be the conductor of $\varphi_{0}$. We can find a positive integer $m$ such that if $\zeta$ is a root of unity in $K$ and $\zeta-1 \in \mathfrak{r a}^{m}$, then $\zeta=1$; see [S97, Lemma 24.3 (1)]. Put $U=\left\{b \in K_{\mathbf{a}}^{\times} \prod_{v \in \mathbf{h}} \mathfrak{r}_{v}^{\times} \mid b-1 \prec \mathfrak{r a}{ }^{m}\right\}$. For $x=a b c$ with $a \in K^{\times}, b \in U$,
and $c \in F_{\mathbf{A}}^{\times}$, define $\varphi(x)=b_{\mathbf{a}}^{\nu}\left|b_{\mathbf{a}}\right|^{-\nu} \varphi_{0}(c)$. This is a well-defined map of $K^{\times} U F_{\mathbf{A}}^{\times}$ into $\mathbf{T}$. Indeed, suppose $a b c=1$; put $\zeta=a / a^{\rho}$. Then $\zeta=b^{\rho} / b \in U \cap K \subset \mathfrak{r}^{\times}$. Since $\left|\zeta_{v}\right|=1$ for every $v \in \mathbf{a}, \zeta$ is a root of unity. By our choice of $\mathfrak{a}^{m}$, we have $\zeta=1$, so that $a \in F^{\times}$and $b \in F_{\mathbf{A}}^{\times} \cap U$. Thus $\varphi_{0}(c)^{-1}=\varphi_{0}(a b)=\varphi_{0}(b)=b_{\mathbf{a}}^{\nu}\left|b_{\mathbf{a}}\right|^{-\nu}$. This shows that $\varphi$ is a well-defined character of $K^{\times} U F_{\mathbf{A}}^{\times}$. Since $K^{\times} U F_{\mathbf{A}}^{\times}$is a subgroup of $K_{\mathbf{A}}^{\times}$of finite index, by [S97, Lemma 11.15] we can extend $\varphi$ to a T-valued character of $K_{\mathbf{A}}^{\times}$, which is clearly a Hecke character with the desired properties.

A5.2. Lemma. Let $E^{*}=\left\{a \in G L_{n}(K)_{\mathbf{a}} \prod_{v \in \mathbf{h}} G L_{n}\left(\mathfrak{r}_{v}\right) \mid a-1 \prec \mathfrak{r a}\right\}$ and $T=$ $\left\{\operatorname{diag}[a, \widehat{a}] \mid a \in E^{*}\right\}$ with an integral $\mathfrak{g}$-ideal $\mathfrak{a}$, and let $C$ be an open subgroup of $D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]$ containing $T$ and such that $a_{x}-1 \prec \mathfrak{r a}$ for every $x \in C$. Suppose $\mathfrak{a} \supset \mathfrak{c}$ and $\mathfrak{a}^{-1} \mathfrak{c}$ is divisible by the relative discriminant of $K$ over $F$. Then $C=$ $T\left(C \cap\left(G_{1}\right)_{\mathbf{A}}\right)$.

Proof. Given $g \in C$, put $u=\operatorname{det}(g)$. Then $u u^{\rho}=1$ and $u_{v} \in \mathfrak{r}_{v}^{\times}$for every $v \in \mathbf{h}$. Suppose we can find $b \in K_{\mathbf{a}}^{\times}$such that $u=b / b^{\rho}, b-1 \prec \mathfrak{r a}$, and $b_{v} \in \mathfrak{r}_{v}^{\times}$ for every $v \in \mathbf{h}$. Then putting $h=\operatorname{diag}[e, \hat{e}]$ with $e=\operatorname{diag}\left[b, 1_{n-1}\right]$, we see that $h \in T$ and $h^{-1} g \in C \cap\left(G_{1}\right)_{\mathbf{A}}$, which proves our lemma. Clearly the problem is to find $b_{v}$ with the required properties for each $v \in \mathbf{v}$. There is no problem for $v \in \mathbf{a}$. If $v \in \mathbf{h}$ and $v \nmid c$, then $v$ is unramified in $K$, and hence the desired $b_{v}$ exists by virtue of [S97, Lemma 5.11 (1)]. Suppose $v \mid c ;$ put $f=\operatorname{det}\left(a_{g}\right)_{v}$. Since $a_{g}\left(d_{g}\right)^{*}-1 \prec \mathfrak{r c}$ and $u-\operatorname{det}\left(a_{g} d_{g}\right) \prec \mathfrak{r c}$, we see that $f \in \mathfrak{r}_{v}^{\times}, f-1 \in \mathfrak{r}_{v} \mathfrak{a}_{v}$, and $u_{v}-f / f^{\rho} \in \mathfrak{r}_{v} \mathfrak{c}_{v}$. Put $w=u_{v} f^{\rho} / f$. Then $w w^{\rho}=1$ and $w-1 \in \mathfrak{r}_{v} \mathfrak{c}_{v}$. By our asumption on $\mathfrak{a}^{-1} \mathfrak{c}$ and [S97, Lemma 17.5], there exists an element $c \in \mathfrak{r}_{v}^{\times}$such that $c-1 \in \mathfrak{r}_{v} \mathfrak{a}_{v}$ and $w=c / c^{\rho}$. Put $b_{v}=c f$. Then $u_{v}=b_{v} / b_{v}^{\rho}$ and $b_{v}-1 \in \mathfrak{r}_{v} \mathfrak{a}_{v}$ as desired. This completes the proof.

A5.3. In [S97, Section A7] we treated theta series of a hermitian form, and proved transformation formulas for them analogous to Theorems A3.3, A3.8 and Propositions A3.13, A3.17, and A3.19. However, the formulas in [S97] were given in terms of $\left(G_{1}\right)_{\mathbf{A}}$ for $G_{1}=G \cap S L_{2 n}(K)$ with $G=U\left(\eta_{n}\right)$ in Case UT. Let us now show that we can formulate the results in terms of $G_{\mathbf{A}}$. (The group $G$ of [S97, Section A7] is the present $G_{1}$.)

Let the notation be as in [S97, Theorem A7.4]; in particular we recall that with $V=K_{n}^{q}$ we defined the action of $\left(G_{1}\right)_{\mathbf{A}}$ on $\mathcal{S}\left(V_{\mathbf{h}}\right)$, that depends on a hermitian element $H$ in $K_{q}^{q}$. We let $\varepsilon$ denote (instead of $\chi$ we used in [S97, Theorem A7.4]) the quadratic Hecke character of $F$ corresponding to $K / F$. By Lemma A5.1 there exists a Hecke character $\varphi$ of $K$ such that $\varphi=\varepsilon$ on $F_{\mathbf{A}}^{\times}, \varphi(y)=y^{-\mathbf{a}}|y|^{\mathbf{a}}$ for $y \in K_{\mathbf{a}}^{\times}$, and the conductor of $\varphi$ divides a power of the conductor of $\varepsilon$. We are going to show that we can extend the action of $\left(G_{1}\right)_{\mathbf{A}}$ on $\mathcal{S}\left(V_{\mathbf{h}}\right)$ to that of $G_{\mathbf{A}}$. We define $r, s \in \mathbf{Z}^{\mathbf{a}}$ so that $H_{v}$ has signature ( $r_{v}, s_{v}$ ) for every $v \in \mathbf{a}$ and put

$$
\begin{equation*}
J_{H}(\alpha, z)=\prod_{v \in \mathbf{a}} j_{v}(\alpha, z)^{r_{v}}{\overline{j_{v \rho}(\alpha, z)}}^{s_{v}} \quad\left(\alpha \in G_{\mathbf{A}}, z \in \mathcal{H}\right) \tag{A5.1}
\end{equation*}
$$

using the notation of (5.3). We also define $A=\left(A_{v}\right)_{v \in \mathbf{a}}$ and $f\left(z ; u, u^{\prime} ; \lambda\right)$ as in [S97, (A7.3.2)] for $z \in \mathcal{H},\left(u, u^{\prime}\right) \in\left(C_{q}^{n}\right)^{\mathbf{a}} \times\left(C_{q}^{n}\right)$ and $\lambda \in \mathcal{S}\left(V_{\mathbf{h}}\right)$. We defined the action of $G_{1}$ on $\mathcal{H} \times\left(C_{q}^{n}\right)^{\mathbf{a}} \times\left(C_{q}^{n}\right)$ by [S97, (A7.3.4), (A7.3.5)]. Clearly this action can be extended to the action of $G$ by the same formulas. (The symbols $\varphi, \mu\left(\alpha_{\tau}, z_{\tau}\right), \kappa\left(\alpha_{\tau}, z_{\tau}\right)$ there correspond to $\mathbf{a}, \mu\left(\alpha_{v}, z_{v}\right), \lambda\left(\alpha_{v}, z_{v}\right)$ here.)

A5.4. Theorem. Every element $\sigma$ of $G_{\mathbf{A}}$ gives a $\mathbf{C}$-linear automorphism of $\mathcal{S}\left(V_{\mathbf{h}}\right)$, written $\lambda \mapsto{ }^{\sigma} \lambda$ for $\lambda \in \mathcal{S}\left(V_{\mathbf{h}}\right)$, with the following properties:
(0) If $\sigma \in\left(G_{1}\right)_{\mathbf{A}}$, then this action is the same as that of [S97, Theorem A7.4]. (Notice that $G$ there is $G_{1}$ here.)
(1) $f\left(\alpha\left(z ; u, u^{\prime}\right) ;{ }^{\alpha} \lambda\right)=J_{H}(\alpha, z) f\left(z ; u, u^{\prime} ; \lambda\right) \quad$ for every $\quad \alpha \in G$.
(2) ${ }^{\left(\sigma^{\prime} \sigma\right)} \lambda=\sigma^{\prime}\left({ }^{\sigma} \lambda\right)$.
(3) The map $\lambda \mapsto{ }^{\sigma} \lambda$ depends only on $\sigma_{\mathrm{h}}, H$, and the choice of $\varphi$; it does not depend on $\left\{A_{v}\right\}_{v \in \mathbf{a}}$ or $\sigma_{\mathbf{a}}$.
(4) $\left\{\left.\sigma \in G_{\mathbf{A}}\right|^{\sigma}{ }^{\sigma}=\lambda\right\}$ is an open subgroup of $G_{\mathbf{A}}$ for every $\lambda \in \mathcal{S}\left(V_{\mathbf{h}}\right)$.
(5) If $\sigma \in P_{\mathbf{h}}$, then

$$
\left({ }^{\sigma} \lambda\right)(\xi)=\left|\operatorname{det}\left(a_{\sigma}\right)\right|_{K}^{q / 2} \varphi_{\mathbf{h}}\left(\operatorname{det}\left(a_{\sigma}\right)\right)^{q} \mathbf{e}_{\mathbf{h}}\left(2^{-1} \operatorname{tr}\left(\xi^{*} H \xi a_{\sigma} b_{\sigma}^{*}\right)\right) \lambda\left(\xi a_{\sigma}\right)
$$

(6) If $\eta=\left[\begin{array}{cc}0 & -1_{n} \\ 1_{n} & 0\end{array}\right]$, we have

$$
\left({ }^{\eta} \lambda\right)(x)=i^{p}\left|N_{F / \mathbf{Q}}\left(\operatorname{det}\left(2 H^{-1}\right)\right)\right|^{n} \int_{V_{\mathbf{h}}} \lambda(y) \mathbf{e}_{\mathbf{h}}\left(-2^{-1} \operatorname{Tr}_{K / F}\left(\operatorname{tr}\left(y^{*} H x\right)\right)\right) d y
$$

Here $p=n \sum_{v \in \mathbf{a}}\left(r_{v}-s_{v}\right)$ and $d y$ is the Haar measure on $V_{\mathbf{h}}$ such that the measure of $\prod_{v \in \mathbf{h}}\left(\mathfrak{r}_{v}\right)_{n}^{q}$ is $\left|D_{K}\right|^{-n q / 2}, D_{K}$ being the discriminant of $K$.

Proof. Assertions (1~6) for $\sigma \in\left(G_{1}\right)_{\mathbf{A}}$ are those given in [S97, Theorem A7.4]. In order to distinguish them from the present ones, let us denote those old ones for $\sigma \in\left(G_{1}\right)_{\mathbf{A}}$ by $(1)_{1},(2)_{1}, \ldots,(6)_{1}$. To define the action of $G_{\mathbf{A}}$ on $\mathcal{S}\left(V_{\mathbf{h}}\right)$, put $Q=\left\{\operatorname{diag}[a, \widehat{a}] \mid a \in G L_{n}(K)\right\}$. For $\lambda \in \mathcal{S}\left(V_{\mathbf{h}}\right)$ and $p=\operatorname{diag}[a, \widehat{a}] \in Q_{\mathbf{A}}$ with $a \in G L_{n}(K)_{\mathbf{A}}$, we define ${ }^{p} \lambda$ by

$$
\begin{equation*}
\left({ }^{p} \lambda\right)(x)=\left|\operatorname{det}\left(a_{\mathbf{h}}\right)\right|_{K}^{q / 2} \varphi_{\mathbf{h}}(\operatorname{det}(a))^{q} \lambda(x a) \quad\left(x \in V_{\mathbf{A}}\right) \tag{A5.2}
\end{equation*}
$$

Clearly ${ }^{p p^{\prime}} \lambda={ }^{p}\left(p^{\prime} \lambda\right)$ for $p, p^{\prime} \in Q_{\mathbf{A}}$. We have $G_{\mathbf{A}}=Q_{\mathbf{A}}\left(G_{1}\right)_{\mathbf{A}}$ and ${ }^{\sigma} \lambda$ is meaningful for $\sigma \in\left(G_{1}\right)_{\mathbf{A}}$. Thus, for $\alpha=p \sigma \in G_{\mathbf{A}}$ with $p \in Q_{\mathbf{A}}$ and $\sigma \in\left(G_{1}\right)_{\mathbf{A}}$ we naturally put ${ }^{\alpha} \lambda={ }^{p}\left({ }^{\sigma} \lambda\right)$. To show that this is well-defined, put $p \sigma=p^{\prime} \sigma^{\prime}$ with $p^{\prime} \in Q_{\mathbf{A}}$ and $\sigma^{\prime} \in\left(G_{1}\right)_{\mathbf{A}}$; put also $h=p^{-1} p^{\prime}$. Then $h=\sigma\left(\sigma^{\prime}\right)^{-1} \in Q \cap\left(G_{1}\right)_{\mathbf{A}}$. Since the definition of ${ }^{p} \lambda$ in (A5.2) is consistent with that of (5) ${ }_{1}$, we have ${ }^{p^{\prime}}\left(\sigma^{\prime} \lambda\right)=$ ${ }^{p}\left({ }^{h}\left(\sigma^{\prime} \lambda\right)\right)={ }^{p}\left(h \sigma^{\prime} \lambda\right)={ }^{p}\left({ }^{\sigma} \lambda\right)$. Thus ${ }^{\alpha} \lambda$ is well-defined, and clearly (3) holds.

Before proving (2), we first make the following observation: Given $w \in\left(G_{1}\right)_{\mathbf{A}}$, put $y=p^{-1} w p$ with $p=\operatorname{diag}[a, \widehat{a}]$. Then $y \in\left(G_{1}\right)_{\mathbf{A}}$. We have to show that ${ }^{p}\left({ }^{y} \lambda\right)={ }^{w}\left({ }^{p} \lambda\right)$, which is equivalent to
(*) If $\lambda^{\prime}(x)=\lambda(x a)$, then $\left({ }^{w} \lambda^{\prime}\right)(x)=\left({ }^{y} \lambda\right)(x a)$, where $x \in V_{\mathbf{A}}$.
In view of $(4)_{1},\left(^{*}\right)$ is true for $w$ in a sufficiently small open subgroup $D$ of $\left(G_{1}\right)_{\mathbf{A}}$ depending on $\lambda$ and $a$. Now $\left(G_{1}\right)_{\mathbf{A}}=G_{1} D$, and $G_{1}$ is generated by $\eta_{n}$ and $P_{1}=\left\{g \in G_{1} \mid c_{g}=0\right\}$. Therefore, in order to prove $\left(^{*}\right)$ for that particular $\lambda$, it is sufficient to prove (*) for $w \in P_{1}$ and $w=\eta_{n}$, and for an arbitrary $\lambda$. If $w \in P_{1}$, $\left.{ }^{*}\right)$ can be derived from (5) ${ }_{1}$ by a straightforward calculation. If $w=\eta_{n}$, we obtain the desired fact from (5) $1_{1}$ and $(6)_{1}$ by observing that $p^{-1} \eta_{n} p=\eta_{n} \operatorname{diag}[t, \widehat{t}]$ with $t=a^{*} a$, and $\varepsilon(\operatorname{det}(t))=1$.

Now, to prove (2), let $\alpha=p \sigma$ as above and let $\beta=r \tau$ with $r=\operatorname{diag}[b, \widehat{b}] \in Q_{\mathbf{A}}$ and $\tau \in\left(G_{1}\right)_{\mathbf{A}}$. Then $\beta \alpha=r p \xi \sigma$ with $\xi=p^{-1} \tau p$. We have shown that ${ }^{p}\left(\xi^{\xi}\right)=$ ${ }^{\tau}\left({ }^{p} \lambda\right)$, from which we immediately obtain ${ }^{(\beta \alpha)} \lambda={ }^{\beta}\left({ }^{\alpha} \lambda\right)$. As for (1), we have it for $\alpha \in G_{1}$; if $\alpha \in Q$, we can verify it by a direct calculation. Since $G=Q G_{1}$, we
obtain (1) in the general case in view of (2). Similarly (5) follows from (5) ${ }_{1}$ and (A5.2), since $P=Q P_{1}$. To prove (4), define $E^{*}$ and $T$ as in Lemma A5.2 with an integral ideal $\mathfrak{a}$ and put $C=\{\alpha \in D[\mathfrak{a}, \mathfrak{a}] \mid \alpha-1 \prec \mathfrak{r a}\}$. (We are taking $\mathfrak{c}$ in that lemma to be $\mathfrak{a}^{2}$ here.) Given $\lambda \in \mathcal{S}\left(V_{\mathbf{h}}\right)$, we can take $\mathfrak{a}$ so that $\lambda(x a)=\lambda(x)$ for every $a \in E^{*}$ and ${ }^{\sigma} \lambda=\lambda$ for every $\sigma \in C \cap\left(G_{1}\right)_{\mathbf{A}}$. We also assume that $\mathfrak{a}$ is divisible by the relative discriminant of $K$ over $F$ and by the conductor of $\varphi$. Then ${ }^{p} \lambda=\lambda$ for every $p \in T$ by (A5.2). Since $C=T\left(C \cap\left(G_{1}\right)_{\mathbf{A}}\right)$, we see that ${ }^{\alpha} \lambda=\lambda$ for every $\alpha \in C$, which proves (4). This completes the proof, as (6) is (6) $)_{1}$.

A5.5. We now consider the case $q=n$; thus, hereafter $V=K_{n}^{n}$. We take $\mu \in \mathbf{Z}^{\mathbf{b}}$ such that $\mu_{v} \geq 0$ for every $v \in \mathbf{b}$ and $\mu_{v} \mu_{v \rho}=0$ for every $v \in \mathbf{a}$, and take also $\tau \in S^{+}$. We then put $l=\mu+n \mathbf{a}$ and

$$
\begin{equation*}
\theta(z, \lambda)=\sum_{\xi \in V} \lambda(\xi) \operatorname{det}(\xi)^{\mu \rho} \mathbf{e}_{\mathbf{a}}^{n}\left(\xi^{*} \tau \xi z\right) \quad\left(\lambda \in \mathcal{S}\left(V_{\mathbf{h}}\right), z \in \mathcal{H}\right) \tag{A5.3}
\end{equation*}
$$

Here we understand that $\operatorname{det}(\xi)^{\mu \rho}=1$ for every $\xi$ if $\mu=0$. Taking $\tau$ to be $H$ in the above theorem and fixing a Hecke character $\varphi$, we have an action of $G_{\mathbf{A}}$ on $\mathcal{S}\left(V_{\mathbf{h}}\right)$. Now we have

$$
\begin{equation*}
\theta\left(\alpha z,{ }^{\alpha} \lambda\right)=j_{\alpha}^{l}(z) \theta(z, \lambda) \text { for every } \alpha \in G \tag{A5.4}
\end{equation*}
$$

Indeed, $\theta(z, \lambda)$ is the series of [S97, (A7.13.1)], and (A5.4) for $\alpha \in G_{1}$ was given in [S97, (A7.13.4)]. If $\alpha=\operatorname{diag}[a, \widehat{a}]$ with $a \in G L_{n}(K)$, then (A5.4) follows immediately from (A5.2), so that (A5.4) holds in general. From (4) of the above theorem we see that $\theta(z, \lambda)$ belongs to $\mathcal{M}_{l}$. We have, for the same reason as in (A3.16),
(A5.5) $\theta(z, \lambda)$ is a cusp form if $\mu \neq 0$.
We now define a function $\theta_{\mathrm{A}}$ on $G_{\mathbf{A}}$ by

$$
\begin{equation*}
\theta_{\mathbf{A}}(x, \lambda)=j_{x}^{l}(\mathbf{i})^{-1} \theta\left(x(\mathbf{i}),{ }^{x} \lambda\right) \quad\left(x \in G_{\mathbf{A}}\right) \tag{A5.6}
\end{equation*}
$$

Then from (A5.4) we easily obtain

$$
\begin{equation*}
\theta_{\mathbf{A}}(\alpha x w, \lambda)=j_{w}^{l}(\mathbf{i})^{-1} \theta_{\mathbf{A}}\left(x,{ }^{w} \lambda\right) \text { if } \alpha \in G, w \in G_{\mathbf{A}}, \text { and } w(\mathbf{i})=\mathbf{i} . \tag{A5.7}
\end{equation*}
$$

Let $\omega$ be a Hecke character of $K$, and $\mathfrak{f}$ the conductor of $\omega$. Taking an integral $\mathfrak{g}$-ideal $\mathfrak{a}$, put

$$
\begin{equation*}
\omega^{\prime}=\omega \varphi^{-n}, \quad \mathfrak{h}=\mathfrak{f} \cap \mathfrak{a}, \quad R^{*}=\left\{w \in V_{\mathbf{A}} \mid w-1 \prec \mathfrak{r a}\right\} . \tag{A5.8a}
\end{equation*}
$$

We then define $\lambda_{0} \in \mathcal{S}\left(V_{\mathbf{h}}\right)$ as follows: $\lambda_{0}(x)=\omega_{\mathrm{h}}\left(\operatorname{det}(x)^{-1}\right)$ if $x \in R^{*}$ and $x_{v} \in G L_{n}\left(\mathfrak{r}_{v}\right)$ for every $v \mid \mathfrak{h} ; \quad \lambda_{0}(x)=0$ otherwise. Fixing $r \in G L_{n}(K)_{\mathbf{h}}$, put $\lambda_{1}(x)=\omega\left(\operatorname{det}(r)^{-1}\right) \lambda_{0}\left(r^{-1} x\right)$ for $x \in V_{\mathrm{h}}$. In [S97, Proposition A7.16] and its proof, we found a group

$$
\begin{equation*}
C^{\prime}=\left\{w \in D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right] \mid a_{w}-1 \prec \mathfrak{r a}\right\} \tag{A5.8b}
\end{equation*}
$$

with a fractional $\mathfrak{g}$-ideal $\mathfrak{b}$ and an integral $\mathfrak{g}$-ideal $\mathfrak{c}$ divisible by $\mathfrak{h}$ and the conductor of $\varphi$ such that

$$
\begin{equation*}
{ }^{w^{w}} \lambda_{1}=\omega_{\mathrm{c}}^{\prime}\left(\operatorname{det}\left(a_{w}\right)\right)^{-1} \lambda_{1} \text { for every } w \in\left(G_{1}\right)_{\mathbf{A}} \cap C^{\prime} \tag{A5.9}
\end{equation*}
$$

(Corrections to [S97, Proposition A7.16]: $D\left[\mathfrak{b}^{-1}, \mathfrak{b c}\right]$ there should be replaced by its subgroup defined by $a_{w}-1 \prec \mathfrak{a} ; \omega_{0}^{\prime}\left(\operatorname{det}\left(a_{\gamma}\right)\right)$ in [S97, (A7.16.1)] should be $\omega_{\mathfrak{c}}^{\prime}\left(\operatorname{det}\left(a_{w}\right)\right)$. Here we are taking $\mathfrak{r a}$ as $\mathfrak{a}$ there.) Now (A5.9) is true for every $w \in C^{\prime}$ if $\mathfrak{c}$ is suitably chosen. Indeed, by Lemma A5.2 we can shoose $\mathfrak{c}$ so that $C^{\prime}=T\left(C^{\prime} \cap\left(G_{1}\right)_{\mathbf{A}}\right)$. Since (A5.9) for $w \in T$ follows from (A5.2), we have (A5.9) for every $w \in C^{\prime}$.

Put $\mathbf{g}(x)=\theta_{\mathbf{A}}\left(x, \lambda_{1}\right)$. Then (A5.7) and the last fact show that $\mathbf{g} \in \mathcal{M}_{l}\left(C^{\prime}, \omega^{\prime}\right)$ with the notation of $\S 20.1$, since ( 20.3 b ) can easily be verified by means of (A5.6). Let $p=\operatorname{diag}[q, \widehat{q}]$ with $q \in G L_{n}(K)_{\mathbf{h}}$ and let $g_{p}$ be the $p$-component of $\mathbf{g}$ in the sense of (20.3b). We are going to show that

$$
\begin{equation*}
g_{p}(z)=\omega^{\prime}(\operatorname{det}(q))^{-1}|\operatorname{det}(q)|_{K}^{n / 2}, \sum_{\xi \in V \cap r R^{*} q^{-1}} \omega_{\mathbf{a}}(\operatorname{det}(\xi)) \omega^{*}\left(\operatorname{det}\left(r^{-1} \xi q\right) \mathfrak{r}\right) \operatorname{det}(\xi)^{\mu \rho} \mathbf{e}_{\mathbf{a}}^{n}\left(\xi^{*} \tau \xi z\right) . \tag{A5.10}
\end{equation*}
$$

Here it is understood that $\omega_{\mathbf{a}}(b) \omega^{*}(b \mathfrak{r})$ for $b=0$ and a fractional ideal $\mathfrak{x}$ denotes $\omega^{*}(\mathfrak{x})$ or 0 according as $\mathfrak{f}=\mathfrak{r}$ or $\mathfrak{f} \neq \mathfrak{r}$. Indeed, for $z=y(\mathbf{i})$ with $y \in G_{\mathrm{a}}$ we have $j_{y}^{l}(\mathbf{i})^{-1} g_{p}(z)=\mathbf{g}(p y)=j_{y}^{l}(\mathbf{i})^{-1} \theta\left(z,{ }^{p} \lambda_{1}\right)$ by (A5.6), so that $g_{p}(z)=\theta\left(z,{ }^{p} \lambda_{1}\right)$. Therefore, by (A5.2) and (A5.3) we obtain

$$
g_{p}(z)=|\operatorname{det}(q)|_{K}^{n / 2} \varphi(\operatorname{det}(q))^{n} \sum_{\xi \in V} \lambda_{1}(\xi q) \operatorname{det}(\xi)^{\mu \rho} \mathbf{e}_{\mathbf{a}}^{n}\left(\xi^{*} \tau \xi z\right)
$$

If $\operatorname{det}(\xi) \neq 0$, then $\lambda_{1}(\xi q) \neq 0$ only when $\xi \in r R^{*} q^{-1}$ and $\operatorname{det}\left(r^{-1} \xi q\right) \mathfrak{r}$ is prime to $\mathfrak{h}$, in which case

$$
\begin{aligned}
\omega(\operatorname{det}(q)) \lambda_{1}(\xi q) & =\omega_{\mathbf{h}}\left(\operatorname{det}\left(r^{-1} \xi q\right)\right) \omega_{\mathbf{a}}(\operatorname{det}(\xi)) \omega_{\mathfrak{h}}\left(\operatorname{det}\left(r^{-1} \xi q\right)\right)^{-1} \\
& =\omega_{\mathbf{a}}(\operatorname{det}(\xi)) \omega^{*}\left(\operatorname{det}\left(r^{-1} \xi q\right) \mathfrak{r}\right)
\end{aligned}
$$

Checking also the case of $\xi$ with zero determinant, we obtain (A5.10). On the other hand, taking (f,r) of (20.9f) to be ( $\mathbf{g}, q$ ) here, we find that

$$
\begin{align*}
c_{\mathbf{g}}(\sigma, q)= & |\operatorname{det}(q)|_{K}^{n / 2} \omega^{\prime}(\operatorname{det}(q))^{-1}  \tag{A5.11}\\
& \cdot \sum_{\xi} \omega_{\mathbf{a}}(\operatorname{det}(\xi)) \omega^{*}\left(\operatorname{det}\left(r^{-1} \xi q\right) \mathfrak{r}\right) \operatorname{det}(\xi)^{\mu \rho}
\end{align*}
$$

where $\xi$ runs over $V \cap r R^{*} q^{-1}$ under the condition that $\xi^{*} \tau \xi=\sigma$.
Now in the setting of $\S 22.3$, take $\left(\chi^{-1}, \mathfrak{e}\right)$ there to be $(\omega, \mathfrak{a})$ here. Then we obtain $\mathbf{g}$ in Case UT in that $\S$, since (A5.11) gives exactly (22.15).

In Case SP the matter is simpler. We take $\theta$ of (A3.21) with $\chi^{-1}$ as $\omega$; we also take $(r, q)$ to be ( $p, r$ ) in (A3.23). By Proposition A3.19 this corresponds to an element $\mathbf{g}$ of $\mathcal{M}_{l}\left(C^{\prime}, \psi^{\prime}\right)$ with $\psi^{\prime}=\chi^{-1} \rho_{\tau}$ and a suitable $C^{\prime}$. Combining (A3.23) with (20.9e), we obtain (22.15). Notice that $\operatorname{det}\left(a_{\beta} d_{\beta}\right)-1 \in \mathfrak{c}$ for $\beta$ in (A3.23), since $\operatorname{diag}\left[r^{-1},{ }^{t} r\right] \beta \in C$ with $r$ there.

## A6. Estimate of the Fourier coefficients of a modular form

A6.1. Let us first recall a basic fact on reduction theory of symmetric and hermitian matrices. We consider $S$ of (16.1a) and $S^{+}, S_{\mathrm{a}}^{+}$of (22.1a, b) in Cases SP and UT, and let $R$ denote the group of all upper triangular elements of $G L_{n}(K)$ whose diagonal elements are all equal to 1 . (Here $K$ is as in $\S 3.5$, and so $K=F$ in Case SP.) We embed $F$ naturally into $F_{\mathbf{a}}$ and extend the map $\operatorname{Tr}_{F / \mathbf{Q}}: F \rightarrow \mathbf{Q}$
to an $\mathbf{R}$-linear map of $F_{\mathbf{a}}=\mathbf{R}^{\mathbf{a}}$ into $\mathbf{R}$, and denote it by the same symbol $\operatorname{Tr}_{F / \mathbf{Q}}$. Given a positive number $r>1$, we denote by $\Delta_{r}$ the set of all diagonal matrices $\operatorname{diag}\left[\delta_{1}, \ldots, \delta_{n}\right]$ with $\delta_{i} \in F_{\mathbf{a}}^{\times}$such that
(A6.1) $\quad\left(\delta_{i}\right)_{v}>0, \quad r^{-1} \leq\left(\delta_{i}\right)_{v} / \operatorname{Tr}_{F / \mathbf{Q}}\left(\delta_{i}\right) \leq r, \quad$ and $\quad \operatorname{Tr}_{F / \mathbf{Q}}\left(\delta_{i}\right) \leq r \operatorname{Tr}_{F / \mathbf{Q}}\left(\delta_{i+1}\right)$ for every $i$ and every $v \in \mathbf{a}$. (We of course ignore $\delta_{n+1}$.) For a compact subset $C$ of $R$ we define a Siegel set $\mathfrak{S}$ by

$$
\begin{equation*}
\mathfrak{S}=\mathfrak{S}(r, C)=\left\{\tau^{*} d \tau \mid \tau \in C, d \in \Delta_{r}\right\} \tag{A6.2}
\end{equation*}
$$

Let $U$ be a subgroup of $G L_{n}(\mathfrak{r})$ of finite index. Then we can choose $r, C$, and a finite subset $B$ of $G L_{n}(K) \cap \mathfrak{r}_{n}^{n}$ so that

$$
\begin{equation*}
S_{\mathbf{a}}^{+}=\bigcup_{b \in B} \bigcup_{u \in U} u^{*} \widehat{b} \mathfrak{S}(r, C) b^{-1} u \tag{A6.3}
\end{equation*}
$$

This and (A6.5) below are well-known. To state another basic fact on $\Gamma \backslash \mathcal{H}^{n}$ for any congruence subgroup $\Gamma$ of $G^{n}$, take any nonempty open subset $X$ of $S_{\mathrm{a}}$ and an element $y_{0}$ of $S_{\mathbf{a}}^{+}$, put

$$
\begin{equation*}
T=\left\{x+i y \in \mathcal{H}^{n} \mid x \in X, y_{0}<y \in S_{\mathbf{a}}^{+}\right\} \tag{A6.4}
\end{equation*}
$$

where we write $y_{0}<y$ (and $y>y_{0}$ ) if $y_{v}>\left(y_{0}\right)_{v}$ for every $v \in \mathbf{a}$. Then there exists a finite subset $A$ of $G$ such that

$$
\begin{equation*}
\mathcal{H}^{n}=\bigcup_{\alpha \in A} \Gamma \alpha T \tag{A6.5}
\end{equation*}
$$

A6.2. Lemma. Let $L$ be a $\mathfrak{g}$-lattice in $S$, and $U$ a subgroup of $G L_{n}(\mathfrak{r})$ of finite index. Then there exists a positive constant $M$ with the following property: Given $h \in L \cap S^{+}$, there exists an element $u$ of $U$ such that $\operatorname{tr}\left(\left(u^{*} h u\right)_{v}^{-1}\right) \leq M$ for every $v \in \mathbf{a}$.

Proof. By (A6.3), given $h \in L \cap S^{+}$, we have $b^{*} u^{*} h u b \in \mathfrak{S}(r, C)$ for some $b \in B$ and $u \in U$. Put $b^{*} u^{*} h u b=\tau^{*} d \tau$ with $\tau \in C$ and $d=\operatorname{diag}\left[\delta_{1}, \ldots, \delta_{n}\right] \in$ $\Delta_{r}$. Then $\left(u^{*} h u\right)^{-1}=b \tau^{-1} d^{-1} \widehat{\tau} b^{*}$. From (A6.1) we obtain $\left(\delta_{i}\right)_{v} \leq r \operatorname{Tr}_{F / \mathbf{Q}}\left(\delta_{i}\right) \leq$ $r^{2} \operatorname{Tr}_{F / \mathbf{Q}}\left(\delta_{i+1}\right) \leq r^{3}\left(\delta_{i+1}\right)_{v}$, and hence $\left(\delta_{i}\right)_{v}^{-1} \leq r^{3}\left(\delta_{1}\right)_{v}^{-1}$. Since $b \tau^{-1}$ belongs to a compact set independent of $h$, we have

$$
\begin{equation*}
\operatorname{tr}\left(\left(u^{*} h u\right)_{v}^{-1}\right) \leq M^{\prime}\left(\delta_{1}\right)_{v}^{-1} \tag{*}
\end{equation*}
$$

with a positive constant $M^{\prime}$ independent of $h$. Now $\tau$ is upper triangular, so that $\delta_{1}=\left(b^{*} u^{*} h u b\right)_{11}$ and this belongs to a fractional ideal $\mathfrak{a}$ depending only on $L$ and $U$. Thus $\operatorname{Tr}_{F / \mathbf{Q}}\left(\delta_{1}\right) \in \operatorname{Tr}_{F / \mathbf{Q}}(\mathfrak{a})$, and $\left(\delta_{1}\right)_{v} \geq r^{-1} \operatorname{Tr}_{F / \mathbf{Q}}\left(\delta_{1}\right) \geq M^{\prime \prime}$ with a positive constant $M^{\prime \prime}$ independent of $h$. Combining this with $\left(^{*}\right)$ we obtain our lemma.

A6.3. Lemma. Let $L$ be a $\mathfrak{g}$-lattice in $F^{t}$ and $f\left(x_{1}, \ldots, x_{t}\right)$ be a nonzero polynomial in $t$ indeterminates $x_{1}, \ldots, x_{t}$ with complex coefficients of degree $d_{i}$ with respect to $x_{i}$ for each $i$. Then there exists a positive constant $M$ depending only on $L$ and $\left\{d_{i}\right\}$ with the following property: Given $\xi \in F_{\mathbf{a}}^{t}$, there exists an element $b$ of $L$ such that $f(b) \neq 0$ and $\left|\left(\xi_{i}-b_{i}\right)_{v}\right| \leq M$ for every $i$ and every $v \in \mathbf{a}$.

Proof. We may assume that $L=\mathfrak{b}^{t}$ with a fractional ideal $\mathfrak{b}$. For $0<r \in \mathbf{R}$ and $a \in F_{\mathbf{a}}$ put $B_{r}(a)=\left\{b \in \mathfrak{b}| |(a-b)_{v} \mid \leq r\right.$ for every $\left.v \in \mathbf{a}\right\}$. We can find $r$ such that $\# B_{r}(a)>\operatorname{Max}\left(d_{1}, \ldots, d_{t}\right)$ for every $a \in F_{\mathbf{a}}$. Now given $\xi \in F_{\mathbf{a}}^{t}$, we can find $b_{1} \in B_{r}\left(\xi_{1}\right)$ such that $f\left(b_{1}, x_{2}, \ldots, x_{t}\right) \neq 0$; then we find $b_{2} \in B_{r}\left(\xi_{2}\right)$ such that $f\left(b_{1}, b_{2}, x_{3}, \ldots, x_{t}\right) \neq 0$. Eventually we find $b_{i} \in B_{r}\left(\xi_{i}\right)$ for $1 \leq i \leq t$ such that $f\left(b_{1}, \ldots, b_{t}\right) \neq 0$. This proves our lemma.

A6.4. Proposition. Let $k$ be a weight and let $f(z)=\sum_{h \in S} c(h) \mathbf{e}_{\mathbf{a}}^{n}(h z) \in \mathcal{M}_{k}$ in Cases SP and UT; put $m=k$ in Case SP and $m=\left(k_{v}+k_{v \rho}\right)_{v \in \mathbf{a}}$ in Case UT. Then the following assertions hold:
(1) If $f$ is a cusp form, then $\delta(z)^{m / 2} f(z)$ is bounded on $\mathcal{H}^{n}$, and $|c(h)| \leq$ $M\left|\operatorname{det}(h)^{m / 2}\right|$ for every $h \in S$ with a constant $M$ depending only on $f$.
(2) If $m_{v} \neq m_{v^{\prime}}$ for some $v, v^{\prime} \in \mathbf{a}$, then $f$ is a cusp form. Consequently, if $f$ is a not a cusp form, then $m=\kappa \mathbf{a}$ with $0<\kappa \in 2^{-1} \mathbf{Z}$.
(3) If $m=\kappa \mathbf{a}$ with $0<\kappa \in 2^{-1} \mathbf{Z}$, then $|c(h)| \leq M \operatorname{det}(h)^{\kappa \mathbf{a}}$ for every $h \in S^{+}$ with a constant $M$ depending only on $f$.

Proof. The first half of (1) and (2) were stated in [S97, Proposition 10.6] and proved in [S97, §§A4.9 and A4.10] for integral $k$. The case of half-integral $k$ can be reduced to the case of integral $k$ by considering $f^{2}$. To prove the second half of (1), we first take $M>0$ so that $\left|\delta(z)^{m / 2} f(z)\right| \leq M$ on the whole $\mathcal{H}$. Then $|f(x+i y)| \leq$ $M \operatorname{det}(y)^{-m / 2}$. Therefore, from (5.24) we obtain $|c(h)| \leq M^{\prime} \operatorname{det}(y)^{-m / 2} \mathbf{e}_{\mathbf{a}}^{n}(-i h y)$ with a positive constant $M^{\prime}$. Since $f$ is a cusp form, $c(h) \neq 0$ only if $h \in S^{+}$. Thus taking $y=h^{-1}$, we find the desired estimate of $c(h)$ as stated in (1).

Before proving (3), we make two elementary observations. Let $g(z)=\sum_{h \in S_{+}} a(h)$ $\cdot \mathbf{e}_{\mathbf{a}}^{n}(h z) \in \mathcal{M}_{k}$. The series is absolutely convergent, and so $|g(z)| \leq \sum_{h}\left|a(h) \mathbf{e}_{\mathbf{a}}^{n}(h z)\right|$ $=\sum_{h}\left|a(h) \mathbf{e}_{\mathbf{a}}^{n}(i h y)\right|$. If $y_{0}<y \in S_{\mathbf{a}}^{+}$, then $\operatorname{tr}\left(h\left(y-y_{0}\right)\right)_{v} \geq 0$ for every $v \in \mathbf{a}$, so that $\mathbf{e}_{\mathbf{a}}^{n}\left(i h\left(y-y_{0}\right)\right) \leq 1$. Thus $\mathbf{e}_{\mathbf{a}}^{n}(i h y) \leq \mathbf{e}_{\mathbf{a}}^{n}\left(i h y_{0}\right)$, and hence $|g(z)| \leq \sum_{h}\left|a(h) \mathbf{e}_{\mathbf{a}}^{n}\left(i h y_{0}\right)\right|$. This means that every element of $\mathcal{M}_{k}$ is bounded on the set $T$ of (A6.4). Next we observe that

$$
\begin{equation*}
|\operatorname{det}(x+i y)| \geq \operatorname{det}(y) \quad\left(x \in S_{v}, y \in S_{v}^{+}\right) \tag{A6.6}
\end{equation*}
$$

To show this, put $\varepsilon=y^{-1 / 2}$ and take a unitary matrix $u$ so that $u \varepsilon x \varepsilon u^{*}=$ $\operatorname{diag}\left[d_{1}, \ldots, d_{n}\right]$ with $d_{i} \in \mathbf{R}$. Then $\left|\operatorname{det}\left((x+i y) y^{-1}\right)=\left|\operatorname{det}\left(\varepsilon x \varepsilon+i 1_{n}\right)\right|=\right.$ $\prod_{\nu=1}^{n}\left|d_{\nu}+i\right| \geq 1$, which proves (A6.6).

To prove (3), we first assume that $k$ is integral. Take $\Gamma$ so that $f \in \mathcal{M}_{k}(\Gamma)$; put $D=\bigcup_{\alpha \in A} \Gamma \alpha$ with $A$ as in (A6.5) and $f_{\beta}=f \|_{k} \beta$ for $\beta \in D$. Since $\left\{f_{\beta} \mid \beta \in\right.$ $D\}=\left\{f_{\alpha} \mid \alpha \in A\right\}$, by our observation we have $\left|f_{\beta}(z)\right| \leq M_{1}$ for every $\beta \in$ $D$ and every $z \in T$. Here and in the following we denote by $M_{1}, M_{2}, \ldots$ some positive constants depending only on $\Gamma, k$, and $f$. Given $z \in \mathcal{H}^{n}$, take $\beta \in D$ and $w \in T$ so that $z=\beta w$; put $c=c_{\xi}$ and $d=d_{\xi}$ with $\xi=\beta^{-1}$. Take a $\mathfrak{g}$-lattice $L$ in $S$ so that $\left[\begin{array}{ll}1 & s \\ 0 & 1\end{array}\right] \in \Gamma$ for every $s \in L$. By Lemma A6.3 we can find $M_{2}$ with the property that for every $x \in S_{\mathbf{a}}$ we have $\operatorname{det}(c s+d) \neq 0$ and $\left\|(s-x)_{v}\right\| \leq M_{2}$ with some $s \in L$, where $\|X\|=\operatorname{Max}_{i, j}\left\{\left|X_{i j}\right|\right\}$. Putting $z=x+i y$, we choose such an $s$ for this particular $x$ and put $\zeta=\left[\begin{array}{cc}0 & 1 \\ -1 & s\end{array}\right]$. Then $\zeta \in$ $G$ and $j_{\beta}(w)^{-1}=j_{\beta}\left(\beta^{-1} z\right)^{-1}=j_{\xi}(z)=j\left(\xi \zeta^{-1} \zeta, z\right)=j\left(\xi \zeta^{-1}, \zeta z\right) j(\zeta, z)$, so that $f(z)=j_{\beta}^{k}(w) f_{\mathcal{\beta}}(w)=j^{k}(\zeta, z)^{-1} j^{k}\left(\xi \zeta^{-1}, \zeta z\right)^{-1} f_{\mathcal{\beta}}(w)$. Put $u=\zeta z$. Then $j\left(\xi \zeta^{-1}, \zeta z\right)=(c s+d) u-c=(c s+d)\left(u-(c s+d)^{-1} c\right), j_{\zeta}(z)=\operatorname{det}(s-z)$, and $u=(s-z)^{-1}$. Since $(c s+d)^{-1} c \in S$, we have, by (A6.6),
$\left(^{*}\right) \quad\left|\operatorname{det}\left(u-(c s+d)^{-1} c\right)\right| \geq \delta(u)=\left|j_{\zeta}(z)\right|^{-2} \delta(z)=|\operatorname{det}(s-z)|^{-2} \delta(z)$.
Therefore, for $k=\kappa \mathbf{a}$ we have

$$
\begin{equation*}
|f(z)| \leq M_{1}|\operatorname{det}(c s+d)|^{-\kappa \mathbf{a}}|\operatorname{det}(z-s)|^{\kappa \mathbf{a}} \delta(z)^{-\kappa \mathbf{a}} \tag{**}
\end{equation*}
$$

Since $\xi \in \alpha^{-1} \Gamma$ with some $\alpha \in A$, we have $c s+d \prec \mathfrak{x}$ with a fractional ideal $\mathfrak{x}$ in $K$ depending only on $A$ and $\Gamma$. Therefore $|\operatorname{det}(c s+d)|^{\kappa \mathbf{a}} \geq M_{3}$. Since $\left\|(x-s)_{v}\right\| \leq M_{2}$, we see that $|\operatorname{det}(x-s+i y)|^{\kappa \mathbf{a}}$ as a polynomial function of $y$ has bounded coefficients. Now $\left|y_{\mu \nu}\right|^{2} \leq\left|y_{\mu \mu} y_{\nu \nu}\right| \leq\left(y_{\mu \mu}+y_{\nu \nu}\right)^{2} / 4 \leq \operatorname{tr}(y)^{2}$, and hence $|\operatorname{det}(x-s+i y)|^{\kappa \mathbf{a}} \leq$ $M_{4}\left(1+\operatorname{tr}(y)^{n}\right)^{\kappa \mathrm{a}}$. Thus from (**) we obtain

$$
\begin{equation*}
|f(x+i y)| \leq M_{5} \operatorname{det}(y)^{-\kappa \mathbf{a}}\left(1+\operatorname{tr}(y)^{n}\right)^{\kappa \mathbf{a}} \tag{A6.7}
\end{equation*}
$$

on the whole $\mathcal{H}^{n}$. We have assumed that $k$ is integral. If $k$ is half-integral, applying (A6.7) to $f^{2}$, we find that (A6.7) is also true for half-integral $k$. Combining (A6.7) with (5.24), we have

$$
\begin{equation*}
|c(h)| \leq M_{6} \mathbf{e}_{\mathbf{a}}^{n}(-i h y) \operatorname{det}(y)^{-\kappa \mathbf{a}}\left(1+\operatorname{tr}(y)^{n}\right)^{\kappa \mathbf{a}} \tag{***}
\end{equation*}
$$

for every $h \in S$. Suppose $c(h) \neq 0$ and $h \in S^{+}$. By (5.21), we have $c\left(u^{*} h u\right)=c(h)$ for every $u$ in a subgroup $U_{1}$ of $G L_{n}(\mathfrak{r})$ of finite index that depends only on $\Gamma$. We can also find a $\mathfrak{g}$-lattice $L_{1}$ in $S$ such that $c(h) \neq 0$ only if $h \in L_{1}$. By Lemma A6.2 we can find $u \in U_{1}$ so that $\operatorname{tr}\left(\left(u^{*} h u\right)_{v}^{-1}\right) \leq M_{7}$. This means that to prove (3), we may assume that $\operatorname{tr}\left(h_{v}^{-1}\right) \leq M_{7}$. Take $y=h^{-1}$ in ( ${ }^{* * *)}$. Then $|c(h)| \leq M_{8} \operatorname{det}(h)^{\kappa \mathrm{a}}$. This completes the proof of (3).

A6.5. Lemma. For $\mathfrak{S}=\mathfrak{S}(r, C)$ as in (A6.2) the following assertions hold:
(1) For every $x \in \mathfrak{S}$ and every $v \in \mathbf{a}$ we have

$$
\operatorname{det}\left(x_{v}\right) \leq\left(x_{11} \cdots x_{n n}\right)_{v} \leq M \cdot \operatorname{det}\left(x_{v}\right)
$$

with a positive constant $M$ depending only on $\mathfrak{S}$.
(2) There exists a positive constant $M^{\prime}$ depending only on $\mathfrak{S}$ with the property that $\left(x_{i i}\right)_{v} /\left(x_{i i}\right)_{v^{\prime}} \leq M^{\prime}$ for every $x \in \mathfrak{S}$, every $i$, and every $v, v^{\prime} \in \mathbf{a}$.
(3) For every $i \leq n$, the number of elements $e \in \mathfrak{g}^{\times}$such that $e x_{i i}=x_{i i}^{\prime}$ for some $x, x^{\prime} \in \mathfrak{S}$ is bounded by a constant depending only on $\mathfrak{S}$.

Proof. Let $x=\tau^{*} d \tau$ with $\tau \in C$ and $d=\operatorname{diag}\left[\delta_{1}, \ldots, \delta_{r}\right] \in \Delta_{r}$. Focusing our attention on one $v \in \mathbf{a}$, let us drop the subscript $v$. We easily see that $x_{11}=\delta_{1}$, and $x_{i i}$ is a linear form of $\delta_{1}, \ldots, \delta_{i}$ with coefficients in a compact set. Therefore from (A6.1) we obtain $x_{i i} \leq M \delta_{i}$ with $M$ depending only on $\mathfrak{S}$. Also we easily see that $x_{i i}-\delta_{i} \geq 0$. Thus $\operatorname{det}(x)=\delta_{1} \cdots \delta_{n} \leq x_{11} \cdots x_{n n} \leq M^{n} \delta_{1} \cdots \delta_{n} \leq M^{n} \operatorname{det}(x)$, which proves (1). As for (2), we have seen that $\left(\delta_{i}\right)_{v} \leq\left(x_{i i}\right)_{v} \leq M\left(\delta_{i}\right)_{v}$. From (A6.1) we see that $\left(\delta_{i}\right)_{v^{\prime}} /\left(\delta_{i}\right)_{v} \leq r^{2}$ for every $v, v^{\prime} \in \mathbf{a}$. Therefore $\left(x_{i i}\right)_{v^{\prime}} /\left(x_{i i}\right)_{v} \leq$ $M\left(\delta_{i}\right)_{v^{\prime}} /\left(\delta_{i}\right)_{v} \leq M r^{2}$, which proves (2). To prove (3), suppose $e x_{i i}=x_{i i}^{\prime}$ for some $x, x^{\prime} \in \mathfrak{S}$ and $e \in \mathfrak{g}^{\times}$. Then $\left(e x_{i i}\right)_{v} /\left(e x_{i i}\right)_{v^{\prime}} \leq M r^{2}$, ad hence $e_{v} / e_{v^{\prime}} \leq M^{2} r^{4}$. Multiplying the inequalities for all $v^{\prime} \in \mathbf{a}$, we find that $M^{-2} r^{-4} \leq e_{v} \leq M^{2} r^{4}$ for every $v \in \mathbf{a}$. Clearly the number of such $e \in \mathfrak{g}^{\times}$is finite. This proves (3).

A6.6. Proposition. Let $U$ be a subgroup of $G L_{n}(\mathfrak{r})$ of finite index and let $S(\mathfrak{a})$ be defined by (16.1b) with a fractional ideal $\mathfrak{a}$ in $K$. Further let $\mathcal{R}=\left[S(\mathfrak{a}) \cap S^{+}\right] / U$ (see §22.1). Then $\sum_{\sigma \in \mathcal{R}} \operatorname{det}(\sigma)^{-s \mathbf{a}}$ is convergent for $\operatorname{Re}(s)>\lambda_{n}$ with $\lambda_{n}$ of §22.1.

Proof. Take $B$ and $\mathfrak{S}=\mathfrak{S}(r, C)$ as in (A6.3). We first prove that for every $\mathfrak{g}$-lattice $\Lambda$ in $S$ the series $\sum_{h \in \mathfrak{G} \cap \Lambda} \operatorname{det}(h)^{-s \mathbf{a}}$ is convergent for $\operatorname{Re}(s)>\lambda_{n}$. For that purpose, take $M$ as in Lemma 6.5 (1). We may assume that $\Lambda \subset S(\mathfrak{a})$ with some $\mathfrak{a}$. Given $h \in \mathfrak{S} \cap \Lambda$, put $h_{i}=h_{i i}$. Then $\operatorname{det}\left(h_{v}\right)^{-1} \leq M\left(h_{1} \cdots h_{n}\right)_{v}^{-1}$ and $\left|\left(h_{i j}\right)_{v}\right|^{2} \leq\left|\left(h_{i} h_{j}\right)_{v}\right|$. Now for fixed positive constants $T, T_{0}$ and $t=\left(t_{v}\right)_{v \in \mathbf{a}} \in \mathbf{R}^{\mathbf{a}}$ such that $t_{v}>T_{0}$ and $t_{v} / t_{v^{\prime}} \leq T$ for every $v, v^{\prime} \in \mathbf{a}$ we have

$$
\begin{equation*}
\#\left\{a \in \mathfrak{a}\left|\left|a_{v}\right| \leq t_{v} \text { for every } v \in \mathbf{a}\right\} \leq M_{1} \prod_{v \in \mathbf{a}} t_{v}^{[K: F]}\right. \tag{A6.8}
\end{equation*}
$$

with a positive constant $M_{1}$ independent of $t$. This will be shown at the end of the proof. Now take $M^{\prime}$ as in Lemma A6.5 (2); put $e=[F: \mathbf{Q}]$. Then $\left(h_{i}\right)_{v}^{e} \geq M^{\prime-e} N_{F / \mathbf{Q}}\left(h_{i}\right) \geq M^{\prime-e} N(\mathfrak{a} \cap F)$, since $h_{i} \in \mathfrak{a} \cap F$. Thus we can take $\left(h_{i} h_{j}\right)_{v}$ as $t_{v}$ in (A6.8). Therefore the number of elements $h \in \mathfrak{S} \cap \Lambda$ with given $h_{1}, \ldots, h_{n}$ as their diagonal elements are at most $M_{1}^{n(n-1) / 2} \cdot N_{F / \mathbf{Q}}\left(h_{1} \cdots h_{n}\right)^{[K: F](n-1) / 2}$. Observe that $[K: F](n-1) / 2=\lambda_{n}-1$. In view of Lemma 6.5 (3), our series in question converges if $\sum_{g}\left|N_{F / \mathbf{Q}}(g)\right|^{\lambda_{n}-1-s}$ is convergent, where $g$ runs over $\mathfrak{a} \cap F^{\times} \bmod -$ ulo multiplication by the elements of $\mathfrak{g}^{\times}$. Since such a series is convergent for $\operatorname{Re}(s)>\lambda_{n}$, we obtain the desird result concerning $\sum_{h \in \mathfrak{S} \cap \Lambda} \operatorname{det}(h)^{-s \mathbf{a}}$. Returning to $\sum_{\sigma \in \mathcal{R}} \operatorname{det}(\sigma)^{-s \mathbf{a}}$ as in our proposition, we may assume, by (A6.3), that $\sigma$ in the sum belongs to $\widehat{b} \subseteq b^{-1}$ for some $b \in B$. Thus it is sufficient to consider the sum of $\operatorname{det}(\sigma)^{-s \mathbf{a}}$ for all $\sigma \in \widehat{b} \mathfrak{S} b^{-1} \cap S(\mathfrak{a})$, or equivalently, the sum for all $\sigma \in \mathfrak{S} \cap b^{*} S(\mathfrak{a}) b$, to which the above result is applicable. Therefore we obtain our proposition.

To prove (A6.8), given $t$, put $\tau=\operatorname{Max}_{v \in \mathbf{a}} t_{v}$. Then $T_{0} \leq \tau \leq T t_{v}$, and so the problem can be reduced to the case in which $t_{v}=\tau$ for every $v \in \mathbf{a}$. Let $\left\{\alpha_{i}\right\}_{i=1}^{d}$ be a $Z$-basis of $\mathfrak{a}$, where $d=[K: \mathbf{Q}]$. Suppose $\left|\left(\sum_{i=1}^{d} m_{i} \alpha_{i}\right)_{v}\right| \leq \tau$ with $m_{i} \in \mathbf{Z}$ for every $v \in \mathbf{a}$. Then $\left|m_{i}\right| \leq A \tau$ with a positive constant $A$ depending only on $\left\{\alpha_{i}\right\}_{i=1}^{d}$. Change $A$ for $\operatorname{Max}\left(A, T_{0}^{-1}\right)$. Then $A \tau \geq 1$, and hence the number of such $m_{i}$ is $\leq 3 A \tau$. Thus we obtain (A6.8).

A6.7. Let us now investigate the convergence of the series of (22.4). As explained at the beginning of $\S 22.3, c_{\mathbf{f}}(\sigma, q)$ (resp $c_{\mathbf{g}}(\sigma, q)$ ) for $\sigma \in S$ are the Fourier coefficients of an element of $\mathcal{S}_{k}\left(\operatorname{resp} \mathcal{M}_{\ell}\right)$. By Proposition A6.4 we have $\left|c_{\mathbf{f}}(\sigma, q)\right| \leq M \operatorname{det}(\sigma)^{m / 2}$ and $\left|c_{\mathbf{g}}(\sigma, q)\right| \leq M^{\prime} \operatorname{det}(\sigma)^{m^{\prime} \varepsilon}$ for every $\sigma \in S^{+}$with constants $M$ and $M^{\prime}$ independent of $\sigma$, where $m$ and $m^{\prime}$ are as in Proposition 22.2; $\varepsilon=1 / 2$ if $\mathbf{g}$ is a cusp form and $\varepsilon=1$ otherwise. Therefore

$$
\left|c_{\mathbf{f}}(\sigma, q) c_{\mathbf{g}}(\sigma, q) \operatorname{det}(\sigma)^{-s \mathbf{a}-h}\right| \leq M M^{\prime}\left|\operatorname{det}(\sigma)^{(c-s) \mathbf{a}}\right|
$$

where $c=0$ if $\mathbf{g}$ is a cusp form, and $c$ is the element of $\mathbf{Q}$ such that $m^{\prime}=2 c \mathbf{a}$ if $\mathbf{g}$ is not a cusp form. Also, $c_{\mathbf{f}}(\sigma, q) c_{\mathbf{g}}(\sigma, q) \neq 0$ only if $\sigma$ belongs to a lattice in $S$ depending on $q$. Therefore, by Proposition A6.6, the series of (22.4) is convergent if $\operatorname{Re}(s)$ is sufficiently large.

## A7. The Mellin transforms of Hilbert modular forms

A7.1. This section concerns the case in which $G=S L_{2}(F)$ and $\mathcal{H}=\mathfrak{H}_{1}^{\mathrm{a}}$ in the setting of Section 5 and $\S 10.6$. We put $\mathfrak{g}_{+}^{\times}=\left\{a \in \mathfrak{g}^{\times} \mid a \gg 0\right\}$.

Let $f(z)=\sum_{h \in F} a(h) \mathbf{e}_{\mathbf{a}}(h z) \in \mathcal{M}_{k}$ with an integral or a half-integral weight $k$. As noted in (5.21) and $\S 6.10$, we can find a subgroup $U_{1}$ of $\mathfrak{g}_{+}^{\times}$of finite index such that $f\left(u^{2} z\right)=u^{-k} f(z)$ for every $u \in U_{1}$, so that $a\left(u^{2} h\right)=u^{k} a(h)$ for every $u \in U_{1}$. We now put

$$
\begin{equation*}
D(s, f)=\left[\mathfrak{g}^{\times}: U\right]^{-1} \sum_{h \in F^{\times} / U} a(h)|h|^{-k / 2-s \mathbf{a}} \quad\left(s \in \mathbf{C}^{\times}\right) \tag{A7.1}
\end{equation*}
$$

where $U$ is a subgroup of $\mathfrak{g}_{+}^{\times}$of finite index such that $a(u h)=u^{k / 2} a(h)$ for every $u \in U$. This is formally well-defined independently of the choice of $U$. By Proposition A6.4, $|a(h)| \leq M|h|^{k / 2+\sigma a}$ for $h \neq 0$ with a constant $M$ independent of $h$, where
$\sigma=\kappa / 2$ if $k=\kappa \mathbf{a}$ with $\kappa \in 2^{-1} \mathbf{Z}$ and $\sigma=0$ otherwise. Therefore the sum is convergent for $\operatorname{Re}(s)>1+\sigma$, and defines a holomorphic function of $s$ there.

To find analytic continuation of $D(s, f)$, put

$$
\begin{equation*}
f_{*}(z)=(-i z)^{-k} f\left(-z^{-1}\right) \tag{A7.2}
\end{equation*}
$$

Observe that $\left(f_{*}\right)_{*}=f, f_{*}=i^{c} f \|_{k} \eta$, and $f_{*} \in \mathcal{M}_{k}$, where $c=\sum_{v \in \mathbf{a}}[k]_{v}$.
A7.2. Theorem. The notation being as above, put $R(s, f)=\Gamma_{k}(s) D(s, f)$ with $\Gamma_{k}(s)=\prod_{v \in \mathbf{a}}(2 \pi)^{-s-k_{v} / 2} \Gamma\left(s+\left(k_{v} / 2\right)\right) ;$ put also $f_{*}(z)=\sum_{h \in F} a_{*}(h) \mathbf{e}_{\mathbf{a}}(h z)$. Then $R(s, f)$ can be continued as a meromorphic function of $s$ to the whole $\mathbf{C}$, and satisfies $R(-s, f)=R\left(s, f_{*}\right)$. Moreover, $R(s, f)$ is entire except when $k=\kappa \mathbf{a}$ with $\kappa \in 2^{-1} \mathbf{Z}$, in which case $R(s, f)$ is holomorphic on $\mathbf{C}$ except for possible simple poles at $s=-\kappa / 2$ and $s=\kappa / 2$ with residues $-a(0) R_{F} / 2$ and $a_{*}(0) R_{F} / 2$, respectively, where $R_{F}$ denotes the regulator of $F$. In particular, $R(s, f)$ is entire if $f$ is a cusp form.

Proof. Put $g=f-a(0), g_{*}=f_{*}-a_{*}(0), \varphi(z)=\sum_{h \in A} a(h) \mathbf{e}_{\mathbf{a}}(h z)$ with a fixed complete set of representatives $A$ for $F^{\times} / U, Y=\left\{y \in \mathbf{R}^{\mathbf{a}} \mid y \gg 0\right\}, Y_{1}=$ $\left\{y \in Y \mid y^{\mathbf{a}} \geq 1\right\}, Y_{2}=\left\{y \in Y \mid y^{\mathbf{a}}<1\right\}$. Then $g(z)=\sum_{u \in U} \sum_{h \in A} a(u h) \mathbf{e}_{\mathbf{a}}(u h z)=$ $\sum_{u \in U} u^{k / 2} \varphi(u z)$, so that $g(i y) y^{k / 2+(s-1) \mathbf{a}}=\sum_{u \in U} \varphi(i u y)(u y)^{k / 2+(s-1) \mathbf{a}}$. Thus we have, at least formally,

$$
\begin{gathered}
\int_{Y / U} g(i y) y^{k / 2+(s-1) \mathbf{a}} d y=\int_{Y} \varphi(i y) y^{k / 2+(s-1) \mathbf{a}} d y \\
=\sum_{h \in A} a(h) \int_{Y} \mathbf{e}_{\mathbf{a}}(i h y) y^{k / 2+(s-1) \mathbf{a}} d y=\sum_{h \in A} a(h) \Gamma_{k}(s)|h|^{-k / 2-s \mathbf{a}}=\left[\mathfrak{g}^{\times}: U\right] R(s, f) .
\end{gathered}
$$

Our formal calculation is valid for $\operatorname{Re}(s)>1+\sigma$, because of the convergence of $\sum_{h \in A}\left|a(h) h^{-k / 2-s \mathbf{a}}\right|$. Now we have

$$
\int_{Y / U} g(i y) y^{k / 2+(s-1) \mathbf{a}} d y=\int_{Y_{1} / U} g(i y) y^{k / 2+(s-1) \mathbf{a}} d y+\int_{Y_{2} / U} g(i y) y^{k / 2+(s-1) \mathbf{a}} d y
$$

All three integrals are convergent for $\operatorname{Re}(s)>1+\sigma$. If $y^{\mathbf{a}} \geq 1$ and $\operatorname{Re}(s)<\operatorname{Re}\left(s^{\prime}\right)$, then $\left|y^{s \mathbf{a}}\right| \leq\left|y^{s^{\prime} \mathbf{a}}\right|$, and hence the integral over $Y_{1} / U$ is convergent for every $s \in \mathbf{C}$. To deal with the integral over $Y_{2} / U$, take $U$ so that $a_{*}(u h)=u^{k / 2} a_{*}(h)$ for every $u \in U$. From (A7.2) we obtain $g\left(i y^{-1}\right)=y^{k} g_{*}(i y)-a(0)+a_{*}(0) y^{k}$, and hence

$$
\begin{aligned}
& \int_{Y_{2} / U} g(i y) y^{k / 2+(s-1) \mathbf{a}} d y=\int_{Y_{1} / U} g\left(i y^{-1}\right) y^{-k / 2-(s+1) \mathbf{a}} d y \\
= & \int_{Y_{1} / U} g_{*}(i y) y^{k / 2-(s+1) \mathbf{a}} d y \\
& -a(0) \int_{Y_{1} / U} y^{-k / 2-(s+1) \mathbf{a}} d y+a_{*}(0) \int_{Y_{1} / U} y^{k / 2-(s+1) \mathbf{a}} d y
\end{aligned}
$$

provided $\operatorname{Re}(s)>1+\sigma$ and the last three integrals are convergent. The integral involving $g_{*}$ is convergent for every $s \in \mathbf{C}$. We have $a(0)=a_{*}(0)=0$ if $f$ is a cusp form. If not, then $k=\kappa \mathbf{a}$ with $\kappa \in 2^{-1} \mathbf{Z}$ by Proposition A6.4 (2). Now we need a formula

$$
\begin{equation*}
\int_{Y_{1} / U} y^{-\alpha \mathbf{a}} d y=(\alpha-1)^{-1} 2^{-1}\left[\mathfrak{g}^{\times}: U\right] R_{F} \quad \text { if } \quad \operatorname{Re}(\alpha)>1 \tag{A7.3}
\end{equation*}
$$

which will be proven at the end of the proof. Thus, putting $M=\left[\mathfrak{g}^{\times}: U\right]^{-1}$, we obtain

$$
\begin{align*}
R(s, f)= & M \int_{Y_{1} / U} g(i y) y^{k / 2+(s-1) \mathbf{a}} d y+M \int_{Y_{1} / U} g_{*}(i y) y^{k / 2-(s+1) \mathbf{a}} d y  \tag{A7.4}\\
& -a(0)\left(R_{F} / 2\right)(s+\kappa / 2)^{-1}+a_{*}(0)\left(R_{F} / 2\right)(s-\kappa / 2)^{-1}
\end{align*}
$$

for sufficiently large $\operatorname{Re}(s)$, where the last two terms occur only if $k=\kappa$. Since the integrals over $Y_{1} / U$ are convergent for every $s$, the right-hand side defines a meromorphic function on the whole $\mathbf{C}$ with poles and residues as described in our theorem. To obtain the functional equation, change $f$ for $f_{*}$. Since $\left(f_{*}\right)_{*}=f$, $R\left(s, f_{*}\right)$ can be obtained by exchanging $(g, a(0))$ for $\left(g_{*}, a_{*}(0)\right)$ in (A7.4). Then we easily see that it coincides with $R\left(-s, f_{*}\right)$. This proves our theorem.

To prove (A7.3), take a set of generators $\left\{\varepsilon_{i}\right\}_{i=1}^{r}$ of $U$, where $r=[F: \mathbf{Q}]-$ 1. Take $r+1$ real variables $\left\{t_{i}\right\}_{i=0}^{r}$ and put $y_{v}=t_{0} \exp \left(\sum_{i=1}^{r} t_{i} \log \left|\varepsilon_{i}\right|_{v}\right)$ for $v \in \mathbf{a}$ with $t_{0}>0$. We easily see that the jacobian of the map $t \mapsto y$ is $\pm(r+1) t_{0}^{r} \operatorname{det}\left(\log \left|\varepsilon_{i}\right|_{v}\right)_{i, v}$, where $v$ is restricted to the set of arbitrarily chosen $r$ elements of $\mathbf{a}$. This quantity equals $\pm(r+1) t_{0}^{r} 2^{-1}\left[\mathfrak{g}^{\times}: U\right] R_{F}$. Therefore

$$
\int_{Y_{1} / U} y^{-\alpha \mathbf{a}} d y=(r+1) t_{0}^{r} 2^{-1}\left[\mathfrak{g}^{\times}: U\right] R_{F} \int_{1}^{\infty} t_{0}^{r-(r+1) \alpha} d t_{0}
$$

which gives (A7.3).
A7.3. Proof of Lemma 18.2. Given $\kappa$ and $t$ as in our lemma, put

$$
f(z)=\theta(z, \kappa)=\sum_{\xi \in F} \kappa(\xi) \xi^{t} \mathbf{e}_{\mathbf{a}}\left(\xi^{2} z / 2\right)=\sum_{h \in F} a(h) \mathbf{e}_{\mathbf{a}}(h z)
$$

This can be obtained by putting $n=1$ and $\tau=1 / 2$ in (A3.12). As explained after (A3.14), $f \in \mathcal{M}_{k}$ with $k=t+\mathbf{a} / 2 ; f$ is a cusp form if $t \neq 0$. Clearly $a(h)=\sum_{\xi^{2}=2 h} \kappa(\xi) \xi^{t}$. Take $U$ of (18.1) so that $U \subset \mathfrak{g}_{+}^{\times}$and put $U^{2}=\left\{u^{2} \mid u \in U\right\}$. Observing that $\bigsqcup_{h \in F^{\times} / U^{2}}\left\{\xi \in F \mid \xi^{2}=2 h\right\}$ gives $F^{\times} / U$, we have

$$
\begin{aligned}
& {\left[\mathfrak{g}^{\times}: U\right] D_{t}(2 s, \kappa)=\sum_{h \in F^{\times} / U^{2}} \sum_{\xi^{2}=2 h} \kappa(\xi) \xi^{t}|\xi|^{-t-2 s \mathbf{a}}} \\
& =\sum_{h \in F^{\times} / U^{2}} a(h)|2 h|^{-t / 2-s \mathbf{a}}=\left[\mathfrak{g}^{\times}: U^{2}\right] 2^{-t / 2-s a} D(s-1 / 4, f) .
\end{aligned}
$$

Consequently

$$
\begin{equation*}
R_{t}(s, \kappa)=2^{r} R((s / 2)-(1 / 4), f) \tag{A7.5}
\end{equation*}
$$

where $r=[F: \mathbf{Q}]-1$. Taking $\eta$ as $\alpha$ in (A3.14) and substituting $\eta z$ for $z$, we obtain $\theta\left(z,{ }^{\eta} \kappa\right)=J(\eta, \eta z) f(\eta z)$, and hence $f_{*}(z)=i^{-|t|} \theta\left(z,{ }^{\eta} \kappa\right)$. Now ${ }^{\eta} \kappa$ can be obtained from Theorem A3.3 (6), where we take $S(x, y)=x y$ for $x, y \in F$. Putting $\kappa_{*}=i^{-|t|} .{ }^{\eta} \kappa$, we find $\kappa_{*}$ as given in Lemma 18.2. Then all the assertions of Lemma 18.2 can be derived from Theorem A7.2 in view of (A7.5).

A7.4. Proof of Theorem 18.16 (1). The notation being as in $\S 18.15$, for $0 \leq$ $m \in \mathbf{Z}^{\mathbf{a}}$ and $\ell \in \mathcal{S}\left(K_{\mathbf{h}}\right)$ put

$$
f(z)=\theta(z, \ell)=\sum_{\alpha \in K} \ell(\alpha) \alpha^{m \rho} \mathbf{e}_{\mathbf{a}}\left(\alpha \alpha^{\rho} z\right)=\sum_{h \in F} a(h) \mathbf{e}_{\mathbf{a}}(h z) \quad(z \in \mathcal{H})
$$

This can be obtained from (A5.3) by putting $n=1, \tau=1, \mu=m$, and $\lambda=\ell$. Since $n=1$, we have $S U(\eta)=S L_{2}(F)$ by Lemma 1.3 (2), and hence $f \in \mathcal{M}_{k}$
with $k=m+\mathbf{a}$, where $\mathcal{M}_{k}$ is defined with respect to $G=S L_{2}(F)$. Clearly $a(h)=$ $\sum_{\alpha \alpha^{\rho}=h} \ell(\alpha) \alpha^{m \rho}$. Take $U$ of (18.18) so that $U \subset \mathfrak{g}_{+}^{\times}$and put $U^{2}=\left\{u^{2} \mid u \in U\right\}$. Observing that $\bigsqcup_{h \in F^{\times} / U^{2}}\left\{\alpha \in K \mid \alpha \alpha^{\rho}=h\right\}$ gives $K^{\times} / U$, we have

$$
\begin{aligned}
& {\left[\mathfrak{r}^{\times}: U\right] L_{m}(s, \ell)=\sum_{h \in F^{\times} / U^{2}} \sum_{\alpha \alpha^{\rho}=h} \ell(\alpha) \alpha^{m \rho}|\alpha|^{-m-2 s \mathbf{a}}} \\
& =\sum_{h \in F^{\times} / U^{2}} a(h)|2 h|^{-m / 2-s \mathbf{a}}=\left[\mathfrak{g}^{\times}: U^{2}\right] D_{k}(s-1 / 2, f) .
\end{aligned}
$$

Thus

$$
\begin{equation*}
L_{m}(s, \ell)=2^{r}\left[\mathfrak{r}^{\times}: \mathfrak{g}^{\times}\right]^{-1} D_{k}(s-1 / 2, f), \tag{A7.6}
\end{equation*}
$$

and therefore (1) of Theorem 18.16 follows immediately from Theorem A7.2, since $f$ is a cup form if $m \neq 0$.

We can actually include the case $m=0$ and prove some results for $L_{m}(s, \ell)$ completely parallel to Lemma 18.2. The precise statement may be left to the reader, as the task is easy and we do not need it in the present book.

## A8. Certain unitarizable representation spaces

A8.1. Given a finite-dimensional vector space $W$ over $\mathbf{C}$, we denote by $\mathfrak{S}(W)$ the symmetric algebra over $W$, and by $\mathfrak{S}_{e}(W)$, for $0 \leq e \in \mathbf{Z}$, its subspace consisting of all the homogeneous elements of degree $e$. Then $S_{e}(W)$ of $\S 12.3$ and $\mathfrak{S}_{e}(W)$ are dual to each other with respect to the pairing

$$
\begin{equation*}
\left\langle\alpha, x_{1} \cdots x_{e}\right\rangle=\alpha_{*}\left(x_{1}, \ldots, x_{e}\right) \quad\left(x_{i} \in W, \alpha \in S_{e}(W)\right) \tag{A8.1}
\end{equation*}
$$

We now take our setting to be that of Sections 12 through 14 . We denote by $\mathcal{K}$ the maximal compact subgroup of $G_{\mathbf{a}}$ given by $\mathcal{K}=\left\{\alpha \in G_{\mathbf{a}} \mid \alpha(\mathbf{o})=\mathbf{o}\right\}$, where $\mathbf{o}=\left(\mathbf{o}_{v}\right)_{v \in \mathbf{a}}$ with $\mathbf{o}_{v}=0$ for Types AB and CB and $\mathbf{o}_{v}=i 1_{n}$ for Types AT and CT. Then we denote by $\mathfrak{g}$ the Lie algebra of $G_{\mathfrak{a}}$, and by $\mathfrak{k}$ the subalgebra of $\mathfrak{g}$ corresponding to $\mathcal{K}$. Since we have a fixed complex structure of $\mathcal{H}=G / \mathcal{K}$, we have a well-known decomposition $\mathfrak{g}_{\mathbf{C}}=\mathfrak{k}_{\mathbf{C}} \oplus \mathfrak{p}_{+} \oplus \mathfrak{p}_{-}$with the following property:

$$
\begin{equation*}
\left[\mathfrak{k}_{\mathbf{C}}, \mathfrak{p}_{ \pm}\right] \subset \mathfrak{p}_{ \pm}, \quad\left[\mathfrak{p}_{+}, \mathfrak{p}_{+}\right]=\left[\mathfrak{p}_{-}, \mathfrak{p}_{-}\right]=\{0\}, \quad\left[\mathfrak{p}_{+}, \mathfrak{p}_{-}\right]=\mathfrak{k}_{\mathbf{C}} \tag{A8.2}
\end{equation*}
$$

We denote by $\mathfrak{U}$ the universal enveloping algebra of $\mathfrak{g}_{\mathbf{C}}$. We let $\mathfrak{g}$ act on the set $C^{\infty}\left(G_{\mathbf{a}}\right)$ of all $C^{\infty}$ functions on $G_{\mathbf{a}}$ by $(Y f)(x)=(d / d t)_{t=0} f(x \cdot \exp (t Y))$ for $Y \in \mathfrak{g}$ and $x \in G_{\mathbf{a}}$, and extend the action to that of $\mathfrak{U}$ on $C^{\infty}\left(G_{\mathbf{a}}\right)$ as usual. More generally, we view every $\mathfrak{g}$-module as a $\mathfrak{U}$-module, and vice versa. We say that a $\mathfrak{U}$-module $\mathcal{Y}$ is unitarizable if it satisfies the following condition:
(A8.3) $\mathcal{Y}$ has a positive definite hermitian form $\{\}:, \mathcal{Y} \times \mathcal{Y} \rightarrow \mathbf{C}$ such that $\{X g, h\}=-\{g, X h\}$ for every $g, h \in \mathcal{Y}$ and every $X \in \mathfrak{g}$.

Let $\mathcal{Y}$ and $\mathcal{Y}^{\prime}$ be two unitarizable $\mathfrak{U}$-modules. Then we define a $\mathfrak{g}$-module structure on $\mathcal{Y} \otimes \mathbf{C} \mathcal{Y}^{\prime}$ as usual by defining $X\left(g \otimes g^{\prime}\right)=(X g) \otimes g^{\prime}+g \otimes X g^{\prime}$ for $g \in \mathcal{Y}, g^{\prime} \in \mathcal{Y}^{\prime}$, and $X \in \mathfrak{g}$. Defining an inner product on $\mathcal{Y} \otimes_{\mathbf{C}} \mathcal{Y}^{\prime}$ by $\left\{g \otimes g^{\prime}, h \otimes h^{\prime}\right\}=\{g, h\}\left\{h, h^{\prime}\right\}$, we easily see that $\mathcal{Y} \otimes \mathrm{C} \mathcal{Y}^{\prime}$ is unitarizable.

A8.2. Given a representation $\{\rho, V\}$ of $\mathcal{K}$, we denote by $C^{\infty}(\rho)$ the set of all elements $f$ of $C^{\infty}\left(G_{\mathbf{a}}, V\right)$ such that $f\left(x k^{-1}\right)=\rho(k) f(x)$ for every $k \in \mathcal{K}$ and $x \in G_{\mathbf{a}}$. Then $X f=-d \rho(X) f$ for every $X \in \mathfrak{k}$ and $f \in C^{\infty}(\rho)$. Let $A d: G_{\mathbf{a}} \rightarrow$ $\operatorname{Aut}(\mathfrak{g})$ denote the adjoint representation of $G_{\mathbf{a}}$. For $g \in G_{\mathbf{a}}$ the action of $\operatorname{Ad}(g)$ can be extended to $\mathfrak{U}$, which we denote also by $\operatorname{Ad}(g)$. In the following we consider exclusively $A d(k)$ for $k \in K$. It is well-known, and in fact, can easily be verified, that

$$
\begin{equation*}
[(A d(k) B) f](x)=\rho(k)(B f)(x k) \text { if } f \in C^{\infty}(\rho), B \in \mathfrak{U}, k \in \mathcal{K}, x \in G_{\mathbf{a}} \tag{A8.4}
\end{equation*}
$$

In view of (A8.2), $\mathfrak{S}\left(\mathfrak{p}_{ \pm}\right)$can be embedded in $\mathfrak{U}$. Then $\mathfrak{S}\left(\mathfrak{p}_{ \pm}\right)$is stable under $\operatorname{Ad}(k)$ for $k \in \mathcal{K}$.

Let us now describe $\mathfrak{p}_{ \pm}$and $\mathcal{K}$ more explicitly, explaining their relation with the symbols $T$ and $K^{c}$ of $\S 12.1$. For the moment, our group is a single factor $G_{v}$, not $G_{\mathbf{a}}$. For each type there exist an injection $\iota_{0}$ of $\mathcal{K}$ into $K^{c}$ and $\mathbf{C}$-linear bijections $\iota_{ \pm}$of $T$ onto $\mathfrak{p}_{ \pm}$with the following properties:

$$
\begin{align*}
& \iota_{0}(k)=M_{k}(\mathbf{o}) \quad \text { for } \quad k \in \mathcal{K},  \tag{A8.5a}\\
& \mathfrak{p}=\left\{\iota_{+}(u)+\iota_{-}(\bar{u}) \mid u \in T\right\},  \tag{A8.5b}\\
& \iota_{ \pm}(u)^{*}=\iota_{\mp}(\bar{u}) \quad \text { for } \quad u \in T
\end{align*}
$$

(A8.5d) $A d(k) \iota_{+}(u)=\iota_{+}\left({ }^{t} a^{-1} u b^{-1}\right)$ and $A d(k) \iota_{-}(u)=\iota_{-}\left(a u^{t} b\right)$ if $\iota_{0}(k)=(a, b)$.
Here $M_{\alpha}(z)$ is defined by (5.1) and (12.4c); we put $(X+i Y)^{*}=X-i Y$ for $X+i Y \in \mathfrak{g}_{\mathrm{C}}$ with $X, Y \in \mathfrak{g}$; the reader is reminded of the convention made in $\S 12.1$ that $K^{c}$ for Type C is identified with the set of all $(a, a) \in G L_{n}(\mathbf{C})^{2}$.

Now the explicit forms of $\iota_{ \pm}$are as follows:
Types AB and CB: $\quad \iota_{+}(u)=\left[\begin{array}{cc}0 & u \\ 0 & 0\end{array}\right], \quad \iota_{-}(u)=\left[\begin{array}{cc}0 & 0 \\ t^{u} & 0\end{array}\right] \quad(u \in T)$,
Types AT and CT: $\quad \iota_{+}(u)=\frac{1}{2} \beta\left[\begin{array}{cc}0 & i u \\ 0 & 0\end{array}\right] \beta^{-1}, \quad \iota_{-}(u)=\frac{1}{2} \beta\left[\begin{array}{cc}0 & 0 \\ -i .^{t} u & 0\end{array}\right] \beta^{-1}$

$$
\left(\beta=\left[\begin{array}{cc}
-i 1_{n} & i 1_{n} \\
1_{n} & 1_{n}
\end{array}\right], \quad u \in T\right)
$$

For these, see [S90, Section 5], where classical groups of other types are also treated.
Given a representation $\{\rho, V\}$ of $K^{c}$, we have a representation $\rho \circ \iota_{0}$ of $\mathcal{K}$; for simplicity, we write also $\rho$ for $\rho \circ \iota_{0}$. In other words, we identify $k$ with $\iota_{0}(k)$ for $k \in \mathcal{K}$. For such a $\{\rho, V\}$ and $h \in C^{\infty}(\mathcal{H}, V)$, we define $h^{\rho} \in C^{\infty}\left(G_{\mathbf{a}}, V\right)$ by $h^{\rho}(x)=\rho\left(M_{x}(\mathbf{o})\right)^{-1} h(x(\mathbf{o}))$ for $x \in G_{\mathbf{a}}$. Then $h^{\rho} \in C^{\infty}(\rho)$, and moreover, $h \mapsto h^{\rho}$ is a C-linear bijection of $C^{\infty}(\mathcal{H}, V)$ onto $C^{\infty}(\rho)$. Now we have

$$
\begin{align*}
& \iota_{+}\left(u_{1}\right) \cdots \iota_{+}\left(u_{e}\right) g^{\rho}=\left(D_{\rho}^{e} g\right)^{\rho \otimes \tau_{v}^{e}}\left(u_{1}, \ldots, u_{e}\right)  \tag{A8.6a}\\
& \iota_{-}\left(u_{1}\right) \cdots \iota_{-}\left(u_{e}\right) g^{\rho}=\left(E^{e} g\right)^{\rho \otimes \sigma_{v}^{e}}\left(u_{1}, \ldots, u_{e}\right)  \tag{A8.6b}\\
& \quad\left(g \in C^{\infty}(\mathcal{H}, V) ; u_{1}, \ldots, u_{e} \in T_{v}\right) .
\end{align*}
$$

For these, see [S90, pp.259-260, Proposition 7.3] and [S94b, Proposition 2.2]. From (A8.6b) we see that $g$ is holomorphic if and only if $Y g^{\rho}=0$ for every $Y \in \mathfrak{p}_{-}$. More generally, $g$ is nearly holomorphic if and only if there exists a positive integer $t$ such that $Y_{1} \cdots Y_{t} g^{\rho}=0$ for every $Y_{1}, \ldots, Y_{t} \in \mathfrak{p}_{-}$. Now we denote by $H(\rho)$ the set of all $f \in C^{\infty}(\rho)$ such that $Y f=0$ for every $Y \in \mathfrak{p}_{-}$. Then $H(\rho)$ consists of the functions of the form $g^{\rho}$ with all holomorphic maps $g$ of $\mathcal{H}$ into $V$.

In the following treatment, we include the case of half-integral weight. Strictly speaking, in such a case we have to formulate everything in terms of suitable covering groups of $G_{\mathbf{a}}$ and $\mathcal{K}$. Since the formulation is obvious, we do not explicitly employ those coverings, and we merely state the results for both integral and halfintegral weights. See [S94b, p.151] for relevant comments on this point.

We insert here an elementary lemma. Let $\mathcal{G}$ be a unimodular Lie group, $L$ its Lie algebra, and $\Gamma$ a closed unimodular subgroup of $\mathcal{G}$. Then $\Gamma \backslash \mathcal{G}$ has a $\mathcal{G}$-invariant measure $\mu$.

A8.3. Lemma. (1) Let $X \in L$. If $\varphi$ is a $\mathbf{C}$-valued $\Gamma$-invariant $C^{1}$ function on $\mathcal{G}$ such that both $\varphi$ and $X \varphi$ belong to $L^{1}(\Gamma \backslash \mathcal{G})$, then $\int_{\Gamma \backslash \mathcal{G}} X \varphi d \mu=0$.
(2) Let $f$ and $h$ be $\mathbf{C}$-valued $\Gamma$-invariant $C^{1}$ functions on $\mathcal{G}$ and let $X \in L$. Then

$$
\int_{\Gamma \backslash \mathcal{G}} X f \cdot h d \mu=-\int_{\Gamma \backslash \mathcal{G}} f \cdot X h d \mu
$$

provided $f h, X f \cdot h$, and $f \cdot X h$ all belong to $L^{1}(\Gamma \backslash \mathcal{G})$.
Proof. Put $c=\int_{\Gamma \backslash \mathcal{G}} X \varphi d \mu$ and $F(g, t)=(X \varphi)(g \cdot \exp (t X))$ for $g \in \mathcal{G}$ and $t \in \mathbf{R}$. Then $\int_{\Gamma \backslash \mathcal{G}} F(g, t) d \mu(g)=c$ for every $t$. Similarly $\int_{\Gamma \backslash \mathcal{G}}|F(g, t)| d \mu(g)=$ $\int_{\Gamma \backslash \mathcal{G}}|X \varphi| d \mu$, and hence $F(g, t)$ is integrable on $[0,1] \times(\Gamma \backslash \mathcal{G})$. Therefore

$$
\begin{aligned}
c & =\int_{0}^{1} \int_{\Gamma \backslash \mathcal{G}} F(g, t) d \mu(g) d t=\int_{\Gamma \backslash \mathcal{G}} \int_{0}^{1} F(g, t) d t d \mu(g) \\
& =\int_{\Gamma \backslash \mathcal{G}}[\varphi(g \cdot \exp (X))-\varphi(g)] d \mu(g)=0,
\end{aligned}
$$

which proves (1). Here we employ the fact that $F(g, t)=(d / d t) \varphi(g \cdot \exp (t X))$, which holds only under the assumption that $\varphi$ is $C^{1}$. Assertion (2) can be obtained by taking $f h$ to be $\varphi$ in (1).

A8.4. Theorem. Let $\rho(a, b)=\operatorname{det}(b)^{\kappa}$ for $(a, b) \in \mathfrak{K}_{0}$ with $\kappa \in 2^{-1} \mathbf{Z}^{\mathbf{a}}$ in Case SP and $\kappa \in \mathbf{Z}^{\mathbf{a}}$ otherwise; let $\kappa_{0}=\operatorname{Min}_{v \in \mathbf{a}} \kappa_{v}$. Suppose that the following condition is satisfied:
(A8.7) $\kappa_{0} \geq n / 2$ in Case $S P, \kappa_{0} \geq n$ in Case $U T$, and $\kappa_{v} \geq \operatorname{Min}\left(m_{v}, n_{v}\right)$ for every $v \in \mathbf{a}$ such that $G_{v}$ is not compact in Case UB.
Then the following assertions hold:
(1) For any nonzero $f \in H(\rho)$ the $\mathfrak{U}$-module structure of $\mathfrak{U} f$ is completely determined by $\kappa$ independently of the choice of $f$. Moreover, $w \mapsto w f$ for $w \in$ $\mathfrak{S}\left(\mathfrak{p}_{+}\right)$gives a bijection of $\mathfrak{S}\left(\mathfrak{p}_{+}\right)$onto $\mathfrak{U} f$.
(2) Such $\mathfrak{U} f$ is irreducible as a $\mathfrak{U}$-module.
(3) Such $\mathfrak{U} f$ is unitarizable.
(4) If $\{$,$\} is an inner product on \mathfrak{U} f$ as in (A8.3), then $\left\{\mathfrak{S}_{p}\left(\mathfrak{p}_{+}\right) f, \mathfrak{S}_{q}\left(\mathfrak{p}_{+}\right) f\right\}=$ 0 for $p \neq q$.

We can show that (1) does not hold without (A8.7); see Proposition A8.9 below for the precise statement.

Proof. For the present $\rho$, we see that $d \rho$ is a $\mathbf{C}$-linear map of $\mathfrak{k}$ into $\mathbf{C}$. Therefore, in view of (A8.2), a well-known procedure shows that given $y \in \mathfrak{U}$ there exists an element $w \in \mathfrak{S}\left(\mathfrak{p}_{+}\right)$such that $y f=w f$ for every $f \in H(\rho)$. Therefore
$\mathfrak{U} f=\mathfrak{S}\left(\mathfrak{p}_{+}\right) f$, and if we can show that $w \mapsto w f$ for $w \in \mathfrak{S}\left(\mathfrak{p}_{+}\right)$is injective for every such nonzero $f$, then we obtain (1). For that purpose, for $h \in \mathfrak{U} f$ and $k \in \mathcal{K}$, define ${ }^{k} h$ by ${ }^{k} h(x)=\rho(k) h(x k)$ for $x \in G_{\mathbf{a}}$. Then ${ }^{k}(B f)=(\operatorname{Ad}(k) B) f$ by (A8.4), and hence $h \mapsto{ }^{k} h$ for $h \in \mathfrak{U} f$ defines a representation of $K$. We can restrict the action to $\mathfrak{S}_{e}\left(\mathfrak{p}_{+}\right) f$. In view of (A8.5d) the decomposition of $\mathfrak{S}_{e}\left(\mathfrak{p}_{+}\right)$into $\mathcal{K}$ irreducible subspaces follows from the decomposition of (a sum of tensor products of) $\left\{\tau_{r}, S_{r}\left(T_{v}\right)\right\}$ described in Theorem 12.7; in particular, each irreducible subspace has multiplicity 1 . Let $V$ be a $\mathcal{K}$-irreducible subspace of $\mathfrak{S}_{e}\left(\mathfrak{p}_{+}\right)$. Then $V f$ is $\{0\}$ or isomorphic to $V$, and $\mathfrak{U} f$ is the direct sum of all such $V f$, since no two different $V$ 's are isomorphic. Therefore it is sufficient to prove that $V f \neq\{0\}$ for every $V$. Let $\mathfrak{g}^{v}$ be the Lie algebra of $G_{v}$ and let $\mathfrak{p}_{+}^{v}=\mathfrak{p}_{+} \cap \mathfrak{g}_{\mathrm{C}}^{v}$. Then $V=\bigotimes_{v \in \mathbf{a}} W_{v}$ and $e=\sum_{v \in \mathbf{a}} e_{v}$ with an irreducible subspace $W_{v}$ of $\mathfrak{S}_{e_{v}}\left(\mathfrak{p}_{+}^{v}\right)$. Therefore it is sufficient to show that $W_{v} f \neq 0$ for every $v \in \mathbf{a}$. Thus we focus our attention on a single $v$; in other words, we may assume $G_{\mathbf{a}}=G_{v}, \mathfrak{p}_{+}=\mathfrak{p}_{+}^{v}, T=T_{v}, S_{e}(T)=S_{e_{v}}\left(T_{v}\right)$, and $V=W_{v}$; so hereafter we drop the subscript $v$. Let us put $\psi^{(e)}=\psi^{\rho \otimes \tau^{e}}$ for $\psi \in$ $C^{\infty}\left(\mathcal{H}, S_{e}(T)\right)$. By (A8.6a), for $g \in C^{\infty}(\mathcal{H})$ and $X_{i}=\iota_{+}\left(u_{i}\right)$ with $u_{1}, \ldots, u_{e} \in T$, we have

$$
\begin{equation*}
X_{1} \cdots X_{e} g^{\rho}=\left(D_{\rho}^{e} g\right)^{(e)}\left(u_{1}, \ldots, u_{e}\right)=\sum_{Z}\left(\varphi_{Z} D_{\rho}^{e} g\right)^{(e)}\left(u_{1}, \ldots, u_{e}\right) \tag{A8.8}
\end{equation*}
$$

where $Z$ runs over all the irreducible subspaces of $S_{e}(T)$. For each $Z$ we take its dual bases $\left\{\zeta_{i}\right\}$ and $\left\{\omega_{i}\right\}$ with respect to (12.8). Then the value of $\varphi_{Z} D_{\rho}^{e} g$, as an element of $S_{e}(T)$, equals $\sum_{i}\left(D_{\rho}^{Z} g\right)\left(\zeta_{i}\right) \cdot \omega_{i}$ in view of (12.22) and (12.23). (Strictly speaking, we should write $\left\{\zeta_{i}^{Z}\right\}$ and $\left\{\omega_{i}^{Z}\right\}$ for these bases, but we suppress the superscript $Z$, merely remembering that they depend on $Z$.) Therefore the last quantity of (A8.8) equals $\sum_{Z} \sum_{i}\left[\left(D_{\rho}^{Z} g\right)\left(\zeta_{i}\right) \omega_{i}\right]^{(e)}\left(u_{1}, \ldots, u_{e}\right)$. From this we obtain

$$
\begin{equation*}
z g^{\rho}=\sum_{Z} \sum_{i}\left\langle\left[\left(D_{\rho}^{Z} g\right)\left(\zeta_{i}\right) \omega_{i}\right]^{(e)}, z\right\rangle \quad \text { for every } \quad z \in \mathfrak{S}_{e}\left(\mathfrak{p}_{+}\right) \tag{A8.9}
\end{equation*}
$$

Here we identify $\mathfrak{S}_{e}\left(\mathfrak{p}_{ \pm}\right)$with $\mathfrak{S}_{e}(T)$ through $\iota_{ \pm}$, so that (A8.1) takes the form $\left\langle\alpha, x_{1} \cdots x_{e}\right\rangle=\alpha_{*}\left(\iota_{ \pm}\left(x_{1}\right), \ldots, \iota_{ \pm}\left(x_{e}\right)\right)$ for $x_{i} \in \mathfrak{p}_{ \pm}$and $\alpha \in S_{e}(T)$. Given $V$ as above and $0 \neq y \in V$, take a $\mathcal{K}$-irreducible subspace $Z$ of $S_{e}(T)$ so that $\langle Z, V\rangle \neq$ 0 ; put $p=\sum_{i}\left(D_{\rho}^{Z} g\right)\left(\zeta_{i}\right) \omega_{i}$ with this particular $Z$. Then $y g^{\rho}=\left\langle p^{(e)}, y\right\rangle$, since $\left\langle Z^{\prime}, V\right\rangle=0$ for $Z^{\prime} \neq Z$. Thus our task is to show that $p \neq 0$ if $g$ is nonzero and holomorphic. Now $D_{\rho}^{Z} g$ for a holomorphic $g$ is a polynomial function of the function $r(z)$ of (13.4a, b) of degree $e$ with holomorphic coefficients. Its highest homogeneous term is of the form $q(r) g$ with $q \in S_{e}(T, Z)$, as can be seen from (13.25); moreover, $q$ is independent of $g$ as noted there. Therefore the highest homogeneous term of $p$ is $g \sum_{i} q(r)\left(\zeta_{i}\right) \omega_{i}$. To find $q(r)\left(\zeta_{i}\right)$, take $s=0$ and $\alpha=1$ in Lemma 13.9. Then we find that $q(r)\left(\zeta_{i}\right)=(-1)^{e} \psi_{Z}(-\kappa) \zeta_{i}(r)$, since $\xi^{-1}=i r$ for Types AT and CT and $\xi^{-1} \bar{z}=-r$ for Types AB and CB. From the formula for $\psi_{Z}$ in Theorem 12.13 we see that the last quantity is nonzero under the condition on $\kappa$ of our lemma, and hence $p \neq 0$, which completes the proof of (1).

To prove (4), we first observe that $\{A g, h\}=\left\{g,-A^{*} h\right\}$ for $g, h \in \mathfrak{U} f$ and $A \in \mathfrak{g}_{\mathrm{C}}$. Therefore (A8.5c), applied to all factors $G_{v}$ of $G_{\mathbf{a}}$, shows that $\left\{\mathfrak{p}_{+} g, h\right\}=$ $\left\{g, \mathfrak{p}_{-} h\right\}$. From (A8.2) we can easily derive that $\mathfrak{S}_{p}\left(\mathfrak{p}_{-}\right) \mathfrak{S}_{q}\left(\mathfrak{p}_{+}\right) f=0$ if $p>q$, and hence $\left\{\mathfrak{S}_{p}\left(\mathfrak{p}_{+}\right) f, \mathfrak{S}_{q}\left(\mathfrak{p}_{+}\right) f\right\} \subset\left\{f, \mathfrak{S}_{p}\left(\mathfrak{p}_{-}\right) \mathfrak{S}_{q}\left(\mathfrak{p}_{+}\right) f\right\}=0$ if $p>q$. This proves (4).

To prove (2) and (3), it is sufficient to consider again the problem on a single $G_{v}$. Indeed, take $f(x)=\prod_{v \in \mathbf{a}} f_{v}\left(x_{v}\right)$ for $x=\left(x_{v}\right)_{v \in \mathbf{a}} \in G_{\mathbf{a}}$ with $x_{v} \in G_{v}$ and
$f_{v} \in H\left(\rho_{v}\right), \rho_{v}\left(a_{v}, b_{v}\right)=\operatorname{det}\left(b_{v}\right)^{\kappa_{v}}$. Then clearly $\mathfrak{U} f=\bigotimes_{v \in \mathbf{a}} \mathfrak{U}\left(\mathfrak{g}_{\mathrm{C}}^{v}\right) f_{v}$, and hence our problem can be reduced to that on $G_{v}$. Therefore we drop the subscript $v$ in the same sense as above. Now, to prove that $\mathfrak{U} f$ is irreducible, it is sufficient to show that for any $z \in \mathfrak{S}\left(\mathfrak{p}_{+}\right), \neq 0$, there exists an element $w$ of $\mathfrak{S}\left(\mathfrak{p}_{-}\right)$such that $w z f=f$. Since $\mathfrak{p}_{-} f=0$, from (A8.2) we easily see that $\mathfrak{S}_{e}\left(\mathfrak{p}_{-}\right) \mathfrak{S}_{i}\left(\mathfrak{p}_{+}\right) f=0$ for $i<e$ and $\mathfrak{S}_{e}\left(\mathfrak{p}_{-}\right) \mathfrak{S}_{e}\left(\mathfrak{p}_{+}\right) f \subset \mathbf{C} f$. Therefore it is sufficient to show that for any $z \in \mathfrak{S}_{e}\left(\mathfrak{p}_{+}\right), \neq 0$, there exists an element $w$ of $\mathfrak{S}_{e}\left(\mathfrak{p}_{-}\right)$such that $w z f \neq 0$, even with a special choice of $f$. Now $z$ can be expressed as a sum of elements $y$, each $y$ being contained in some $V$ as above. Let $f=g^{\rho}$ with a holomorphic $g$. Then we have seen in (A8.9) that $z f=\left\langle h^{(e)}, z\right\rangle$ with an element $h \in \mathcal{N}_{\rho \otimes \tau^{e}}^{e}$ whose highest term is $\alpha(r) g$ with $\alpha \in S_{e}\left(T, S_{e}(T)\right)$. The above computation of $p$ shows that $\alpha(r)=\sum_{Z} c_{Z} \sum_{i} \zeta_{i}(r) \omega_{i}$ with $c_{Z} \in \mathbf{Q}^{\times}$. Taking $E, \rho \otimes \tau^{e}$, and $h$ in place of $D_{\rho}, \rho$, and $g$ in (A8.8) and (A8.9), and employing (A8.6b) instead of (A8.6a), we find that

$$
w h^{(e)}=\sum_{Z^{\prime}} \sum_{j}\left\langle\left[\left(E^{Z^{\prime}} h\right)\left(\zeta_{j}^{\prime}\right) \omega_{j}^{\prime}\right]^{\rho \otimes \tau^{e} \otimes \sigma^{e}}, w\right\rangle \quad \text { for every } \quad w \in \mathfrak{S}_{e}\left(\mathfrak{p}_{-}\right)
$$

where $Z^{\prime}$ runs over all the irreducible subspaces of $S_{e}(T)$, and $\left\{\zeta_{i}^{\prime}\right\}$ and $\left\{\omega_{i}^{\prime}\right\}$ are dual bases of $Z^{\prime}$; we use $Z^{\prime}$, since we will have to consider the double sum $\sum_{Z} \sum_{Z^{\prime}}$. By (13.14b), for every $\gamma \in Z^{\prime}$ we have $(-1)^{e}\left(E^{Z^{\prime}} h\right)(\gamma)=\gamma(\partial / \partial r) h$, which equals $g \cdot \gamma(\partial / \partial r) \alpha(r)$, since $\gamma(\partial / \partial r)$ annihilates the terms of degree $<e$. Taking $\partial / \partial r$ in place of $\mathcal{D}_{v}$ in (12.28), we obtain $\gamma(\partial / \partial r) \alpha(r)=e![\gamma, \alpha]$. Thus $\sum_{Z^{\prime}} \sum_{j}\left(E^{Z^{\prime}} h\right)\left(\zeta_{j}^{\prime}\right) \omega_{j}^{\prime}=(-1)^{e} e!g \sum_{Z^{\prime}} \sum_{j}\left[\zeta_{j}^{\prime}, \alpha\right] \omega_{j}^{\prime}$. Take $g$ to be constant 1 and evaluate our functions at the identity element 1 of $G_{\mathbf{a}}$. By Lemma $12.8(1),\left[Z, Z^{\prime}\right]=$ 0 for $Z \neq Z^{\prime}$, and hence we obtain

$$
(w z f)(1)=\left\langle\left(w h^{(e)}\right)(1), z\right\rangle=(-1)^{e} e!\sum_{Z} c_{Z} \sum_{i, j}\left[\zeta_{j}, \zeta_{i}\right]\left\langle\omega_{j}, w\right\rangle\left\langle\omega_{i}, z\right\rangle
$$

Pick a $\mathcal{K}$-irreducible subspace $V$ of $\mathfrak{S}_{e}\left(\mathfrak{p}_{+}\right)$so that the $V$-component of $z$ is nonzero. Let $Z$ be the $\mathcal{K}$-irreducible subspace of $S_{e}(T)$ such that $\langle Z, V\rangle \neq 0$, and let $V^{\prime}$ be the image of $V$ under $\iota_{-} \circ\left(\iota_{+}\right)^{-1}$. Take $\left\{\omega_{i}\right\}$ for this particular $Z$ so that $\left\langle\omega_{1}, z\right\rangle=1$ and $\left\langle\omega_{i}, z\right\rangle=0$ for $i>1$. Put $\beta=\sum_{j}\left[\zeta_{j}, \zeta_{1}\right] \omega_{j}$; then $\beta \neq 0$ by Lemma 12.8 (2). Take $w \in V^{\prime}$ so that $\langle\beta, w\rangle \neq 0$. Then $(w z f)(1)=(-1)^{e} e!c_{Z}\langle\beta, w\rangle \neq 0$, which completes the proof of (2).

To prove (3), we first consider $\mathfrak{U} \varphi$ with $\varphi=f_{0}^{\rho}$, where $0 \neq f_{0} \in \mathcal{S}_{\kappa}(\Gamma)$ with a congruence subgroup $\Gamma$ of $G$. Then every element of $\mathfrak{U} \varphi$ is left $\Gamma$-invariant. We then put

$$
\begin{equation*}
\{g, h\}=\int_{\Gamma \backslash G_{\mathbf{a}}} \overline{g(x)} h(x) d x \tag{A8.10}
\end{equation*}
$$

for $g, h \in \mathfrak{U} \varphi$. This is convergent, since $f_{0}$ is a cusp form. Then by Lemma A8.3 we easily see that $\mathfrak{U} \varphi$ is unitarizable. We employ this fact to prove (3). In view of (1), it is sufficient to prove our assertion for a special choice of $f$. In Case SP we can take $G_{\mathbf{a}}=G_{v}=S p(n, \mathbf{R})$. Put $\kappa=q / 2$ with $q \geq n$ in Case SP. Take a real quadratic field $F$ and put $G^{\prime}=S p(n, F)$ and $G_{\mathrm{a}}^{\prime}=G_{v} \times G_{v^{\prime}}$ so that $\left\{v, v^{\prime}\right\}$ is the set of archimedean primes of $F$. By Remark A3.11 (III), we can find a nonzero cusp form $\varphi$ on $G^{\prime}$ of weight $(q / 2)\left(v+v^{\prime}\right)+v^{\prime}$. Since $G^{\prime}$ is dense in $G_{\mathbf{a}}^{\prime}$, we can find an element $\gamma$ of $G^{\prime}$ such that $\varphi(\gamma(\mathbf{i})) \neq 0$. Changing $\varphi$ for $\varphi \| \gamma$, we may assume that $\varphi(\mathbf{i}) \neq 0$. Put $f(x)=\varphi^{\sigma}(x, 1)$ for $x \in G_{v}$, where $\sigma\left(a, a^{\prime}\right)=\operatorname{det}(a)^{q / 2} \operatorname{det}\left(a^{\prime}\right)^{1+(q / 2)}$ for $\left(a, a^{\prime}\right) \in G L_{n}(\mathbf{C})^{2}$. We easily see that
$0 \neq f \in H(\rho)$. Let $\mathfrak{g}^{\prime}$ be the Lie algebra of $G_{\mathbf{a}}^{\prime}$. Then $\mathfrak{U}$ can be embedded into $\mathfrak{U}\left(\mathfrak{g}_{\mathrm{C}}^{\prime}\right)$ and $(\alpha f)(x)=\left(\alpha \varphi^{\sigma}\right)(x, 1)$ for $\alpha \in \mathfrak{U}$. Our proof of (1) shows that the map $\alpha f \mapsto \alpha \varphi^{\sigma}$ is bijective. Since $\varphi$ is a cusp form, $\mathfrak{U}\left(\mathfrak{g}_{\mathrm{C}}^{\prime}\right) \varphi^{\sigma}$ is unitarizable as we have shown by means of (A8.10). Therefore, $\mathfrak{U f}$, being $\mathfrak{g}$-isomorphic to a subspace of $\mathfrak{U}\left(\mathfrak{g}_{\mathrm{C}}^{\prime}\right) \varphi^{\sigma}$, must be unitarizable. This proves (3) in Case SP. Case UT can be handled in a similar way, but here we first prove a lemma applicable to both Cases UT and UB.

A8.5. Lemma. Let $\mathcal{G}=S U(m, n), \mathcal{K}=S U(m, n) \bigcap[U(m) \times U(n)]$, and $\rho(a, b)=$ $\operatorname{det}(b)$ for $(a, b) \in \mathcal{K}$. Then there exist a discrete subgroup $\Gamma$ of $\mathcal{G}$ and a nonzero element $f$ of $H(\rho)$ such that $\Gamma \backslash \mathcal{G}$ is compact and $f(\gamma x)=f(x)$ for every $\gamma \in \Gamma$.

Proof. Take a real quadratic field $F$ and a CM-field $K$ such that $[K: F]=2$; put $\mathbf{a}=\left\{v, v^{\prime}\right\}$. We can then find an element $\mathcal{T}$ of $G L_{m+n}(K)$ such that $\mathcal{T}^{*}=-\mathcal{T}$, $i \mathcal{I}_{v}$ has signature ( $m, n$ ), and $i \mathcal{T}_{v^{\prime}}$ is positive definite. (The easiest way to find such a $\mathcal{T}$ is to take it to be diagonal.) With this $\mathcal{T}$ we consider the injection of $S U(\mathcal{T})$ into $S p(d, \mathbf{Q})$ and the embedding $\varepsilon: \mathcal{H} \rightarrow \mathfrak{H}_{d}$ of $\S 11.6$, where $d=2(m+n)$. Put $\psi(z)=\operatorname{det}(\kappa(z))$ for $z \in \mathcal{H}$ with $\kappa$ of (11.7). From (11.11) we obtain

$$
\begin{equation*}
j(\widetilde{\alpha}, \varepsilon(z))=\psi(\alpha z) j_{\alpha}(z)^{2} \psi(z)^{-1} \quad \text { for every } \quad \alpha \in S U(\mathcal{T}) \tag{A8.11}
\end{equation*}
$$

since $\operatorname{det}\left[M_{\alpha}(z)\right]=j_{\alpha}(z)^{2}$, which can be seen from (3.23), (3.24a), and (5.1). Now take a nonzero element $\theta$ of $\mathcal{M}_{\mathrm{a} / 2}$ on $\mathfrak{H}_{d}$ and a congruence subgroup $\Gamma^{\prime}$ of $S p(d, \mathbf{Q})$ so that $\theta^{2} \in \mathcal{M}_{\mathbf{a}}\left(\Gamma^{\prime}\right)$; put $g(z)=\psi(z)^{-1 / 2} \theta(\varepsilon(z))$ for $z \in \mathcal{H}$ with any fixed square root $\psi(z)^{-1 / 2}$ of $\psi(z)^{-1}$, which is meaningful as a holomorphic function on $\mathcal{H}$. Changing $\theta$ for $\theta \|_{\mathbf{a} / 2} \beta$ with a suitable $\beta \in S p(d, \mathbf{Q})$ if necessary, we may assume that $g \neq 0$. Take a congruence subgroup $\Gamma_{1}$ of $S U(\mathcal{T})$ so that $\left\{\widetilde{\alpha} \mid \alpha \in \Gamma_{1}\right\} \subset \Gamma^{\prime}$. From (A8.11) we see that $j_{\gamma}(z)^{-2} g(\gamma z)^{2}=g(z)^{2}$ for every $\gamma \in \Gamma_{1}$. Then $j_{\gamma}(z)^{-1} g(\gamma z)=\chi(\gamma) g(z)$ for every $\gamma \in \Gamma_{1}$ with a character $\chi: \Gamma_{1} \rightarrow$ $\{ \pm 1\}$. We obtain the desired $\Gamma$ and $f$ by putting $\Gamma=\left\{\gamma \in \Gamma_{1} \mid \chi(\gamma)=1\right\}$ and $f=g^{\rho}$. Indeed, $\Gamma \backslash \mathcal{G}$ is compact, since $i \mathcal{T}_{v^{\prime}}$ is definite.

Returning to the proof of Theorem A8.4 (3) in the unitary cases, take $f$ and $\Gamma$ as in Lemma A8.5 with $\mathcal{G}=G_{v}$, and for $g, h \in \mathfrak{U} f^{\kappa}$ define $\{g, h\}$ by (A8.10), which is meaningful, since $\Gamma \backslash G_{v}$ is compact. By Lemma A8.3, $\mathfrak{U} f^{\kappa}$ is unitarizable. This completes the proof of Theorem A8.4.

A8.6. Lemma. Let $f \in H(\rho)$ and $f^{\prime} \in H\left(\rho^{\prime}\right)$ with $\rho(a, b)=\operatorname{det}(b)^{\kappa}$ and $\rho^{\prime}(a, b)=\operatorname{det}(b)^{\kappa^{\prime}}$ with $\kappa, \kappa^{\prime} \in 2^{-1} \mathbf{Z}^{a}$. Suppose that both $\kappa$ and $\kappa^{\prime}$ satisfy condition (A8.7). Then the $\mathbf{C}$-linear span of the products $\alpha f \cdot \beta f^{\prime}$ for all $\alpha, \beta \in \mathfrak{U}$ is unitarizable.

Proof. Let $\mathfrak{U} f \cdot \mathfrak{U} f^{\prime}$ denote the $\mathbf{C}$-linear span of the products in question. We view $\mathfrak{S}\left(\mathfrak{p}_{+}\right)$and $\mathfrak{S}_{e}\left(\mathfrak{p}_{+}\right)$as representation spaces of $\mathcal{K}$ through $A d$. In the proof of Theorem A8.4 we have seen that $g \mapsto{ }^{k} g$ for $g \in \mathfrak{U} f$ and $k \in \mathcal{K}$ defines a representaion of $\mathcal{K}$, which is equivalent to $\left\{A d, \mathfrak{S}\left(\mathfrak{p}_{+}\right)\right\}$. The same is true for $\mathfrak{U} f^{\prime}$. Naturally we use $\rho^{\prime}$ instead of $\rho$ for the action of $\mathcal{K}$ on $\mathfrak{U} f^{\prime}$; similarly we put $\left({ }^{k} \psi\right)(x)=\rho(k) \rho^{\prime}(k) \psi(x k)$ for $k \in K$ and $\psi \in \mathfrak{U} f \cdot \mathfrak{U} f^{\prime}$. Define a C-linear map $\sigma: \mathfrak{U} f \otimes \mathfrak{U} f^{\prime} \rightarrow \mathfrak{U} f \cdot \mathfrak{U} f^{\prime}$ so that $\sigma(g \otimes h)=g h$. We let $\mathcal{K}$ act on $\mathfrak{U} f \otimes \mathfrak{U} f^{\prime}$ by ${ }^{k}(g \otimes h)={ }^{k} g \otimes{ }^{k} h$. Then clearly ${ }^{k} \sigma(\varphi)=\sigma\left({ }^{k} \varphi\right)$ for $\varphi \in \mathfrak{U} f \otimes \mathfrak{U} f^{\prime}$, and $\sigma$ is a $\mathfrak{g}$-homomorphism. Let $\mathfrak{T}_{p}=\sum_{e+e^{\prime}=p} \mathfrak{S}_{e}\left(\mathfrak{p}_{+}\right) f \otimes \mathfrak{S}_{e^{\prime}}\left(\mathfrak{p}_{+}\right) f^{\prime}$. Then $\mathfrak{U} f \cdot \mathfrak{U} f^{\prime}=$ $\sum_{p=0}^{\infty} \sigma\left(\mathfrak{T}_{p}\right)$. By Theorem A8.4 (3), both $\mathfrak{U} f$ and $\mathfrak{U} f^{\prime}$ are unitarizable, and hence
so is $\mathfrak{U} f \otimes \mathfrak{U} f^{\prime} ;$ let $\{$,$\} be a hermitian inner product with the property of (A8.3)$ on that space. Put $M_{p}=\operatorname{Ker}(\sigma) \cap \mathfrak{T}_{p}$ and $L_{p}=\left\{x \in \mathfrak{T}_{p} \mid\left\{x, M_{p}\right\}=0\right\}$. Then $\mathfrak{T}_{p}=L_{p} \oplus M_{p}$, since $\mathfrak{T}_{p}$ is finite-dimensional. By Theorem A8.4 (4), $\mathfrak{T}_{p}$ and $\mathfrak{T}_{q}$ for $p \neq q$ are orthogonal. Observe also that they have no isomorphic $\mathcal{K}$ irreducible subspaces. Now every element of $\operatorname{Ker}(\sigma)$ is contained in $\sum_{p \in P} \mathfrak{T}_{p}$ with a finite set $P$, and hence $\operatorname{Ker}(\sigma)$, being $\mathcal{K}$-stable, is the sum of some $\mathcal{K}$-irreducible subspaces, each of which is contained in $\mathfrak{T}_{p}$ for some $p$. Thus $\operatorname{Ker}(\sigma)=\bigoplus_{p} M_{p}$, and $\mathfrak{U} f \otimes \mathfrak{U} f^{\prime}=\operatorname{Ker}(\sigma) \oplus L$ with $L=\bigoplus_{p} L_{p}=\left\{\varphi \in \mathfrak{U f \otimes U f ^ { \prime } | \{ \varphi , \operatorname { K e r } ( \sigma ) \} = 0 \} \text { , which } , ~ ( \sigma )}\right.$ is a $\mathfrak{g}$-submodule. Now $\mathfrak{U f} \cdot \mathfrak{U} f^{\prime}$ is $\mathfrak{g}$-isomorphic to $L$, and hence is unitarizable.

A8.7. Lemma. Let $W$ be a subfield of $\mathbf{C}$ such that $\Phi \mathbf{Q}_{\mathrm{ab}} \subset W$ in Cases $S P$ and UT and $\overline{\mathbf{Q}} \subset W$ in Case UB, where $\Phi$ is the Galois closure of $K$ in $\mathbf{C}$ over Q. Let $\rho(a, b)=\operatorname{det}(b)^{\kappa}$ for $(a, b) \in \mathfrak{K}_{0}$ with $\kappa \in 2^{-1} \mathbf{Z}^{\mathbf{a}}$ in Case SP and $\kappa \in \mathbf{Z}^{\mathbf{a}}$ otherwise. If $f$ is an element of $\mathcal{N}_{\kappa}^{p}(W)$ such that $\mathfrak{U} f^{\rho}$ is unitarizable, then there exists an element $q$ of $\mathcal{M}_{\kappa}(W)$ such that $\langle f, h\rangle=\langle q, h\rangle$ for every $h \in \mathcal{S}_{\kappa}$.

Proof. Take dual bases $\left\{a_{\nu}\right\}$ and $\left\{b_{\nu}\right\}$ of $T_{v}$ with respect to $(u, v) \mapsto \operatorname{tr}\left({ }^{t} u v\right)$ as in $\S 12.5$ for a fixed $v \in \mathbf{a}$; define an element $\mathcal{L}_{v}$ of $\mathfrak{U}$ by $\mathcal{L}_{v}=\sum_{\nu \in N} \iota_{+}\left(a_{\nu}\right) \iota_{-}\left(b_{\nu}\right)$, which is clearly independent of the bases. We can take these bases so that $b_{\nu}=$ $a_{\nu}=\bar{a}_{\nu}$ for every $\nu$. Take a hermitian form $\{$,$\} on \mathfrak{U} f^{\rho}$ as in (A8.3). Then $\{A g, h\}=\left\{g,-A^{*} h\right\}$ for $g, h \in \mathfrak{U} f^{\rho}$ and $A \in \mathfrak{g}_{\mathbf{C}}$, so that by (A8.5c),

$$
\begin{aligned}
\left\{\mathcal{L}_{v} g, h\right\} & =\left\{\sum_{\nu \in N} \iota_{+}\left(a_{\nu}\right) \iota_{-}\left(a_{\nu}\right) g, h\right\}=-\sum_{\nu \in N}\left\{\iota_{-}\left(a_{\nu}\right) g, \iota_{-}\left(a_{\nu}\right) h\right\} \\
& =\left\{g, \sum_{\nu \in N} \iota_{+}\left(a_{\nu}\right) \iota_{-}\left(a_{\nu}\right) h\right\}=\left\{g, \mathcal{L}_{v} h\right\}
\end{aligned}
$$

Taking in particular $g=h$, we see that if $\mathcal{L}_{v} g=0$, then $\sum_{\nu \in N}\left\{\iota_{-}\left(a_{\nu}\right) g, \iota_{-}\left(a_{\nu}\right) g\right\}$ $=0$, and hence $\iota_{-}\left(a_{\nu}\right) g=0$ for every $\nu$. Thus we obtain the direct part of

$$
\begin{equation*}
\mathcal{L}_{v} g=0 \text { for every } v \in \mathbf{a} \Longleftrightarrow X g=0 \text { for every } X \in \mathfrak{p}_{-} . \tag{A8.12}
\end{equation*}
$$

The converse part is obvious. Let $L_{v}$ denote the operator $L_{\omega, v}$ of (15.3) with $\omega(x, y)=\operatorname{det}(y)^{\kappa}$. Then, for a function $\varphi$ on $\mathcal{H}$, we have

$$
\begin{aligned}
\mathcal{L}_{v} \varphi^{\rho} & =\sum_{\nu \in N} \iota_{+}\left(a_{\nu}\right) \iota_{-}\left(b_{\nu}\right) \varphi^{\rho}=\sum_{\nu \in N} \iota_{+}\left(a_{\nu}\right)\left(E_{v} \varphi\right)^{\rho \otimes \sigma_{v}}\left(b_{\nu}\right) \\
& =\sum_{\nu \in N}\left(D_{\omega \otimes \sigma_{v}, v} E_{v} \varphi\right)^{\rho \otimes \sigma_{v} \otimes \tau_{v}}\left(b_{\nu}, a_{\nu}\right)=-\left(L_{v} \varphi\right)^{\rho}
\end{aligned}
$$

by (A8.6a, b) and (12.30a, b). Given $f$ as in our lemma, take $\Gamma$ so that $f \in \mathcal{N}_{\kappa}^{p}(\Gamma)$ and put $\mathcal{X}=\sum_{v \in \mathbf{a}^{\prime}} \sum_{i=0}^{\infty} \mathbf{C} L_{v}^{i} f$, where $\mathbf{a}^{\prime}=\left\{v \in \mathbf{a} \mid G_{v}\right.$ is not compact $\}$. Then $\mathcal{X}$ is finite-dimensional, since it is a subspace of $\mathcal{N}_{\kappa}^{p}(\Gamma)$. Now $\mathcal{X}^{\rho}=\sum_{v \in \mathbf{a}^{\prime}} \sum_{i=0}^{\infty} \mathbf{C} \mathcal{L}_{v}^{i} f^{\rho}$ $\subset \mathfrak{U} f^{\rho}$, and hence we can diagonalize the $\mathcal{L}_{v}$ on $\mathcal{X}^{\rho}$ simultaneously. Consequently we can diagonalize the $L_{v}$ on $\mathcal{X}$ simultaneously. Let $h \in \mathcal{S}_{\kappa}$. Then $L_{v} h=0$, as $h$ is holomorphic. Suppose $L_{v} g=\lambda_{v} g$ with $g \in \mathcal{X}$ and $0 \neq \lambda_{v} \in \mathbf{C}$ for some $v$. Then $0=\left\langle g, L_{v} h\right\rangle=\left\langle L_{v} g, h\right\rangle=\bar{\lambda}_{v}\langle g, h\rangle$ by Theorem 12.15, and hence $\langle g, h\rangle=0$. Let $\mathcal{Y}=\left\{g \in \mathcal{X} \mid L_{v} g=0\right.$ for every $\left.v \in \mathbf{a}^{\prime}\right\}, \mathcal{X}(W)=\mathcal{N}_{\kappa}^{p}(\Gamma, W) \cap \mathcal{X}$, and $\mathcal{Y}(W)=\mathcal{Y} \cap \mathcal{X}(W)$. By Theorems 14.9 (2) and 14.12 (3), $L_{v}$ maps $\mathcal{X}(W)$ into itself, so that we have a Jordan decomposition $\mathcal{X}(W)=\mathcal{Y}(W) \oplus \mathcal{Z}$ with a subspace $\mathcal{Z}$ over $W$ such that $\mathcal{Z} \otimes_{W} \mathbf{C}$ is spanned by the eigenfunctions of the $L_{v}$ not belonging to $\mathcal{Y}$. Given $W$-rational $f$, let $q$ be the projection of $f$ to $\mathcal{Y}(W)$
with respect to this decomposition. Then $\langle f, h\rangle=\langle q, h\rangle$ for every $h \in \mathcal{S}_{\kappa}$. Now $\mathcal{L}_{v} q^{\rho}=-\left(L_{v} q\right)^{\rho}=0$ for every $v \in \mathbf{a}$, so that $q^{\rho} \in H(\rho)$ by (A8.12), that is, $q$ is holomorphic. Thus $q$ is the desired element of $\mathcal{M}_{\kappa}(W)$.

A8.8. Proof of Lemma 15.8. The notation being as in the lemma, put $f=g^{\rho}$ with $\rho(a, b)=\operatorname{det}(b)^{l}, f^{\prime}=h^{\rho^{\prime}}, \quad s=\Delta_{l^{\prime}}^{p} h$, and $\varepsilon(a, b)=\operatorname{det}(b)^{l^{\prime}+2 p}$. Take $Z$ and $\psi$ as in our lemma, and take $\zeta_{1}=\psi$ in (A8.9); take also $y \in V$ so that $\left\langle\omega_{1}, y\right\rangle=1$. (Notice that $\operatorname{dim}(Z)=\operatorname{dim}(V)=1$.) Then (A8.9) shows that $y f^{\prime}=s^{\varepsilon}$. Thus $(g s)^{\rho \varepsilon}=f s^{\varepsilon} \in \mathfrak{U} f \cdot \mathfrak{U} f^{\prime}$, so that $\mathfrak{U}\left((g s)^{\rho \varepsilon}\right) \subset \mathfrak{U} f \cdot \mathfrak{U} f^{\prime}$. By Lemma A8.6, $\mathfrak{U} f \cdot \mathfrak{U} f^{\prime}$ is unitarizable, so that $\mathfrak{U}\left((g s)^{\rho \varepsilon}\right)$ is unitarizable. Therefore we obtain the desired element $q$ by Lemma A8.7.

A8.9. Proposition. Let $G_{\mathbf{a}}=S p(n, \mathbf{R})$ or $S U(m, n)$; let $\rho(a, b)=\operatorname{det}(b)^{\kappa}$ for $(a, b) \in \mathfrak{K}_{0}$ with $\kappa \in 2^{-1} \mathbf{Z}, 0 \leq \kappa \leq(n-1) / 2$ if $G_{\mathbf{a}}=S p(n, \mathbf{R})$ and $\kappa \in \mathbf{Z}, 0 \leq$ $\kappa \leq \operatorname{Min}\{m, n\}-1$ if $G_{\mathbf{a}}=\operatorname{SU}(m, n)$. Then there exist two nonzero elements $f_{1}$ and $f_{2}$ of $H(\rho)$ such that $\mathfrak{U} f_{1}$ and $\mathfrak{U} f_{2}$ are not $\mathfrak{U}$-isomorphic.

Proof. Put $\nu=2 \kappa+1$ if $G_{\mathbf{a}}=S p(n, \mathbf{R})$ and $\nu=\kappa+1$ if $G_{\mathbf{a}}=S U(m, n)$; put also $\psi(u)=\operatorname{det}_{\nu}(u)$ for $u \in T$. Let $Z$ be the irreducible subspace of $S_{\nu}(T)$ with $\psi$ as its highest weight vector in the sense of Theorem 12.7. By [S94b, (4.9c)], $\left(D_{\rho}^{Z} g\right)(\zeta)=\zeta(\mathcal{D}) g$ with $\mathcal{D}$ of $(12.25)$ for every $\zeta \in Z$ and every $g \in C^{\infty}(\mathcal{H})$. (The symbols $\lambda$ and $h$ there correspond to $\kappa$ and $\nu$ here.) Take $V \subset \mathfrak{S}_{e}\left(\mathfrak{p}_{+}\right)$ corresponding to $Z$ as in the proof of Theorem A8.4. Then, by (A8.9),

$$
y g^{\rho}=\sum_{i}\left\langle\left[\zeta_{i}(\mathcal{D}) g \cdot \omega_{i}\right]^{\rho \otimes \tau^{e}}, y\right\rangle \quad \text { for every } \quad y \in V
$$

Let $f_{1}=g_{1}^{\rho}$ with a nonzero constant $g_{1}$ on $\mathcal{H}$. Then $\zeta_{i}(\mathcal{D}) g_{1}=0$, so that $V f_{1}=0$. Next take $f_{2}=g_{2}^{\rho}$ with $g_{2}(z)=\psi(z)$ for $z \in \mathcal{H}$. Then $\zeta_{i}(\mathcal{D}) g_{2}=\nu!\left[\zeta_{i}, \psi\right]$ by (12.28), and hence, $\left(y f_{2}\right)(1)=\nu!\left\langle\sum_{i}\left[\zeta_{i}, \psi\right] \omega_{i}, y\right\rangle=\nu!\langle\psi, y\rangle$, which is not zero for some $y \in V$. Thus $\mathfrak{U} f_{2}$ is not isomorphic to $\mathfrak{U} f_{1}$, which proves our proposition.

A8.10. We have been dealing with a set of objects $\{G, \mathcal{K}, \mathcal{H}, \mathfrak{g}, \mathfrak{U}\}$. Suppose we have two more sets $\left\{G_{i}, \mathcal{K}_{i}, \mathcal{H}_{i}, \mathfrak{g}_{i}, \mathfrak{U}_{i}\right\}$ of the same type for $i=1,2$; suppose also that there exist an injective homomorphism $I$ of $\left(G_{1} \times G_{2}\right)_{\mathbf{a}}$ into $G_{\mathbf{a}}$ and a holomorphic injection $J$ of $\mathcal{H}_{1} \times \mathcal{H}_{2}$ into $\mathcal{H}$ such that $I\left(\mathcal{K}_{1} \times \mathcal{K}_{2}\right) \subset \mathcal{K}$ and $J(\beta z, \gamma w)=I(\beta, \gamma) J(z, w)$ for $(\beta, \gamma) \in\left(G_{1} \times G_{2}\right)_{\mathbf{a}}$ and $(z, w) \in \mathcal{H}_{1} \times \mathcal{H}_{2}$. To avoid possible confusions, we assume our setting to be one of the following types:

Cases SP and UT: $G_{1}=G^{n}, G_{2}=G^{r}, G=G^{n+r}, \mathcal{H}_{1}=\mathcal{H}^{n}, \mathcal{H}_{2}=\mathcal{H}^{r}, \mathcal{H}=$ $\mathcal{H}^{n+r}, J(z, w)=\operatorname{diag}[z, w]$, and $I(\beta, \gamma)=\beta \times \gamma$ with the notation of Sections 24 and 25 .

Case UB: $G_{1}=G^{\psi}, G_{2}=G^{\varphi}, G=G^{\eta}, \mathcal{H}_{1}=\mathfrak{Z}^{\psi}, \mathcal{H}_{2}=\mathfrak{Z}^{\varphi}$, and $\mathcal{H}=\mathcal{H}_{q+n}^{\mathbf{a}}$ with the notation of Section 26 and [S97, Section 21]; $I$ and $J$ will be described later.

In these cases the map extended to the complexification of $\mathcal{K}_{1} \times \mathcal{K}_{2}$ is a collection of several maps each of which is equivalent to the diagonal embedding of $G L_{r}(\mathbf{C}) \times$ $G L_{s}(\mathbf{C})$ into $G L_{r+s}(\mathbf{C})$ for some ( $\left.r, s\right)$.

We consider $(\rho, \kappa)$ for $G$ as in Theorem A8.4, and define similarly $\rho_{i}$ for $G_{i}$ with the same $\kappa$. Now $\mathfrak{g}_{1} \times \mathfrak{g}_{2}$ can be embedded in $\mathfrak{g}$, so that $\mathfrak{U}_{1} \otimes \mathfrak{U}_{2}$ can be viewed as a subalgebra of $\mathfrak{U}$. For a function $h$ on $G_{\mathbf{a}}$ we define a function $h^{\circ}$ on $\left(G_{1} \times G_{2}\right)_{\mathbf{a}}$ by $h^{\circ}(x, y)=h(I(x, y))$ for $(x, y) \in\left(G_{1} \times G_{2}\right)_{\mathbf{a}}$. Clearly $(\alpha \otimes \beta)\left(h^{\circ}\right)=$
$((\alpha \otimes \beta) h)^{\circ}$ for $(\alpha, \beta) \in \mathfrak{U}_{1} \times \mathfrak{U}_{2}$. We define ${ }^{k} h$, as before, by $\left({ }^{k} h\right)(z)=\rho(k) h(z k)$ for $k \in \mathcal{K}$ and $z \in G_{\mathbf{a}}$, and similarly put $\left({ }^{\left(k, k^{\prime}\right)} g\right)(x, y)=\rho_{1}(k) \rho_{2}\left(k^{\prime}\right) g\left(x k, y k^{\prime}\right)$ for $g \in C^{\infty}\left(\left(G_{1} \times G_{2}\right)_{\mathbf{a}}\right),(x, y) \in\left(G_{1} \times G_{2}\right)_{\mathbf{a}}$, and $\left(k, k^{\prime}\right) \in \mathcal{K}_{1} \times \mathcal{K}_{2}$. Clearly $\left({ }^{I\left(k, k^{\prime}\right)} h\right)^{\circ}={ }^{\left(k, k^{\prime}\right)}\left(h^{\circ}\right)$.

A8.11. Lemma. For every $f \in H(\rho), \neq 0$, the $\left(\mathfrak{g}_{1} \times \mathfrak{g}_{2}\right)$-module $\left\{h^{\circ} \mid h \in \mathfrak{U} f\right\}$ is unitarizable, provided (A8.7) is satisfied for the group $G$.

Proof. Put $N=\left\{g \in \mathfrak{U} f \mid g^{\circ}=0\right\}, T_{p}=\mathfrak{S}_{p}\left(\mathfrak{p}_{+}\right) f, N_{p}=N \cap T_{p}$. By Theorem A8.4 (3), $\mathfrak{U} f$ is unitarizable; let $\{$,$\} be a hermitian inner product on \mathfrak{U} f$ with the property of (A8.3). Put $M_{p}=\left\{x \in T_{p} \mid\left\{x, N_{p}\right\}=0\right\}$. Then $T_{p}=L_{p} \oplus M_{p}$, since $\mathfrak{T}_{p}$ is finite-dimensional. By Theorem A8.4 (4), $T_{p}$ and $T_{q}$ for $p \neq q$ are orthogonal; also they have no isomorphic $\left(\mathcal{K}_{1} \times \mathcal{K}_{2}\right)$-irreducible subspaces. Now every element of $N$ is contained in $\sum_{p \in P} T_{p}$ with a finite set $P$, and hence $N$, being ( $\mathcal{K}_{1} \times \mathcal{K}_{2}$ )-stable, is the sum of some $\left(\mathcal{K}_{1} \times \mathcal{K}_{2}\right)$-irreducible subspaces, each of which is contained in $T_{p}$ for some $p$. Thus $N=\bigoplus_{p} N_{p}$, and $\mathfrak{U} f=M \oplus N$ with $M=\bigoplus_{p} M_{p}=\{g \in \mathfrak{U} f \mid\{g, N\}=0\}$, which is a $\left(\mathfrak{g}_{1} \times \mathfrak{g}_{2}\right)$-submodule. Now $(\mathfrak{U} f)^{\circ}$ is $\left(\mathfrak{g}_{1} \times \mathfrak{g}_{2}\right)$-isomorphic to $M$, and hence is unitarizable.

A8.12. Proof of Lemma 28.7. Let $\rho_{0}(a, b)=\operatorname{det}(b)^{\nu \mathbf{a}}, \rho(a, b)=\operatorname{det}(b)^{2 p+\nu \mathbf{a}}$, and $\rho^{\prime}(a, b)=\operatorname{det}(b)^{m}$. For the same reason as in $\S A 8.8$, we have $\left(\Delta_{\nu \mathbf{a}}^{p} R\right)^{\rho}=\alpha R^{\rho_{0}}$ with an element $\alpha$ of $\mathfrak{U}$. Assuming that $n=r$, recall that $D_{e, e^{\prime}}=B_{e} C_{e^{\prime}}, B_{e} g=$ $\left(D_{\rho}^{Z} g\right)(\varphi)$, and $C_{e^{\prime}} g=\left(E^{W} g\right)\left(\varphi^{\prime}\right)$, where $\varphi$ is as in (25.3) and $\varphi^{\prime}$ is defined by the same formula with $e^{\prime}$ in place of $e$. Taking $\Delta_{\nu \mathbf{a}}^{p} R$ and the present $D_{\rho}^{Z}$ as $g$ and $D_{\rho}^{Z}$ of (A8.9), put $q=\sum_{i}\left(D_{\rho}^{Z} g\right)\left(\zeta_{i}\right) \omega_{i}$ with $\left\{\zeta_{i}\right\}$ and $\left\{\omega_{i}\right\}$ as in that formula. Then $\beta g^{\rho}=\left\langle q^{\rho \otimes \tau^{z}}, \beta\right\rangle$ for every $\beta$ in the $\mathcal{K}$-irreducible subspace $V$ of $\mathfrak{S}_{n|e|}\left(\mathfrak{p}_{+}\right)$such that $\langle Z, V\rangle \neq 0$. Now we view $Z$ as a representation space of the complexification of $\mathcal{K}_{1} \times \mathcal{K}_{2}$. Then $\mathbf{C} \varphi$ is stable under the group action, and so $Z=\mathbf{C} \varphi+Z^{\prime}$ with a subspace $Z^{\prime}$ stable under the group action. Take $\left\{\zeta_{i}\right\}$ so that $\zeta_{1}=\varphi$ and $Z^{\prime}=$ $\sum_{i>1} \mathbf{C} \zeta_{i}$. From (25.3) we obtain $\varphi\left(\operatorname{diag}[a, b] u \cdot \operatorname{diag}\left[a^{\prime}, b^{\prime}\right]\right)=\operatorname{det}\left(b a^{\prime}\right)^{e} \varphi(u)$, and hence from (12.9) we see that $\omega_{1}$ has the same property, and $\sum_{i>1} \mathbf{C} \omega_{i}$ is $\left(\mathcal{K}_{1} \times \mathcal{K}_{2}\right)$ stable. Take an element $\beta$ of $V$ such that $\left\langle\omega_{1}, \beta\right\rangle=1$ and $\left\langle\omega_{i}, \beta\right\rangle=0$ for $i>1$. Now (A8.9) evaluated at $\mathfrak{x} \in G_{\mathbf{a}}^{n+r}$ involves $\tau^{e}(M(\mathfrak{x}, \mathbf{o}))$. However, if we take $\mathfrak{x}=I(x, y)$, from what we said about the $\omega_{i}$ we can easily derive that $\left\langle q^{\rho \otimes \tau^{z}}, \beta\right\rangle$ at $\mathfrak{x}=I(x, y)$ equals $\left(D_{\rho}^{Z} g\right)(\varphi)^{\rho_{1}}(I(x, y))$, where $\rho_{1}(a, b)=\rho(a, b) \operatorname{det}(b)^{e}$. Thus $\left(\beta g^{\rho}\right)^{\circ}=\left[\left(D_{\rho}^{Z} g\right)(\varphi)^{\rho_{1}}\right]^{\circ}$. Taking similarly $E^{W}$ in place of $D_{\rho}^{Z}$, we can find an element $\gamma$ of $\mathfrak{U}$ such that $\left(\gamma \beta g^{\rho}\right)^{\circ}=\left[\left(D_{e, e^{\prime}} g\right)^{\rho^{\prime}}\right]^{\circ}$. (This can be justified, since $D_{\rho}^{Z}$ is considered on $\prod_{e_{v}>0} G_{v}$, and $E^{Z}$ on $\prod_{e_{v}^{\prime}>0} G_{v}$, and $e_{v} e_{v}^{\prime}=0$ for every v.) Thus $\left(\gamma \beta \alpha R^{\rho_{0}}\right)^{\circ}=\left(S_{0}^{\rho^{\prime}}\right)^{\circ}$. We consider $D_{e, e^{\prime}}$ only when $n=r$, and so if $n \neq r$, what we need is merely $\left(\alpha R^{\rho_{0}}\right)^{\circ}=\left(S_{0}^{\rho^{\prime}}\right)^{\circ}$. Now $\pi^{-c} S_{0} \in \mathcal{N}^{p^{\prime}}(\Xi)$ with some $c$ and $p^{\prime}$ by Theorem 14.12 (4), and hence $\pi^{-c} S(z, w)$ as a function of $w$ (resp. $z$ ) belongs to $\mathcal{N}_{m}^{p^{\prime}}(\Gamma)$ (resp. $\mathcal{N}_{m}^{p^{\prime}}\left(\Gamma^{\prime}\right)$ ) with a congruence subgroup $\Gamma$ of $G^{r}$ (resp. $\Gamma^{\prime}$ of $G^{n}$ ). All such nearly holomorphic forms with respect to $\Gamma^{\prime} \times \Gamma$ form a finite-dimensional vector space with a $\Xi$-rational basis as the proof of Lemma 24.11 shows.

Now suppose that $\nu \geq(n+r) / 2$ in Case SP and $\nu \geq n+r$ in Case UT. Then the $\left(\mathfrak{g}_{1} \times \mathfrak{g}_{2}\right)$-module $\left\{h^{\circ} \mid h \in \mathfrak{U} R^{\rho_{0}}\right\}$ is unitarizable by Lemma A8.11. Define $L_{v}$ on $\mathcal{H}^{r}$ and $\mathcal{L}_{v}$ on $G_{\mathbf{a}}^{r}$ as in the proof of Lemma A8.7 with $\kappa=m$. Then we can repeat the proof of Lemma A8.7. To be explicit, put $\mathcal{X}=\sum_{v \in \mathbf{a}} \sum_{i=0}^{\infty} \mathbf{C} L_{v}^{i} S$ and
$\mathcal{X}(\Xi)=\sum_{v \in \mathbf{a}} \sum_{i=0}^{\infty} \Xi L_{v}^{i} S^{\prime}$, where $S^{\prime}=\pi^{-c} S$. Then $\mathcal{X}(\Xi)$ is a finite-dimensional vector space over $\Xi$. For a function $h \in C^{\infty}\left(\mathcal{H}^{n} \times \mathcal{H}^{r}\right)$ define $h^{\rho^{\prime}} \in C^{\infty}\left(G_{\mathbf{a}}^{n} \times\right.$ $G_{\mathbf{a}}^{r}$ ) in an obvious way by restricting $\rho^{\prime}$ to the complexification of $\mathcal{K}_{1} \times \mathcal{K}_{2}$. Then, as in the proof of Lemma A8.7 we have $\mathcal{L}_{v} h^{\rho^{\prime}}=-\left(L_{v} h\right)^{\rho^{\prime}}$, and so $\mathcal{X}^{\rho^{\prime}}=$ $\sum_{v \in \mathbf{a}} \sum_{i=0}^{\infty} \mathbf{C} \mathcal{L}_{v}^{i}\left(\gamma \beta \alpha R^{\rho_{0}}\right)^{\circ}$, since $\left(S_{0}^{\rho^{\prime}}\right)^{\circ}=S^{\rho^{\prime}}$. Thus $\mathcal{X}^{\rho^{\prime}}$ is contained in a unitarizable space. Therefore, by the same procedure as in the proof of Lemma A8.7 and keeping the variables on $\mathcal{H}^{n}$ and $G^{n}$ constant, we can find an element $T$ of $\mathcal{X}(\Xi)$ such that $L_{v} T=0$ for every $v \in \mathbf{a}$ and $\left\langle S^{\prime}(z, w)-T(z, w), f(w)\right\rangle=0$ for every $f \in \mathcal{S}_{m}^{r}$. The unitarizability implies (A8.12) in the present case, so that $T(z, w)$ is holomorphic in $w$. Thus $\pi^{c} T$ gives the desired element $T$ of Lemma 28.7.

A8.13. Proof of Lemma 29.3. We assume that $r_{v}>0$ for every $v \in \mathbf{a}$, since we do not need Lemma 29.3 otherwise. The idea of the proof is the same as in $\S A 8.12$. However, the map $\iota: \mathfrak{Z}^{\psi} \times \mathfrak{Z}^{\varphi}$ of [S97, (6.10.2)] is antiholomorphic in the variable $w$ on $\mathfrak{Z}^{\varphi}$, and so we have to change it for a holomorphic one. Thus put $\widetilde{w}=\left(\widetilde{w}_{v}\right)_{v \in \mathbf{a}}$, where

$$
\widetilde{z}=\left[\begin{array}{c}
-\bar{x}  \tag{A8.13}\\
\bar{y}
\end{array}\right] \quad \text { if } \quad z=\left[\begin{array}{l}
x \\
y
\end{array}\right] \in \mathfrak{Z}_{v}^{\varphi} .
$$

Notice that $\widetilde{\mathbf{i}}_{v}=\mathbf{i}_{v}$. Also, put $P=\left(P_{v}\right)_{v \in \mathbf{a}}, Q=\left(Q_{v}\right)_{v \in \mathbf{a}}, R=\left(R_{v}\right)_{v \in \mathbf{a}}$ with $P_{v}=\operatorname{diag}\left[-1_{r_{v}}, 1_{t_{v}+r_{v}}\right], Q_{v}=\operatorname{diag}\left[R_{v}, 1_{r_{v}}\right], R_{v}=\operatorname{diag}\left[1_{r_{v}},-1_{t_{v}}\right]$. For $\alpha=\left(\alpha_{v}\right) \in$ $\prod_{v \in \mathbf{a}} U\left(\varphi_{v}^{\prime}\right)$ with $\varphi_{v}^{\prime}$ as in (26.2), put $\widetilde{\alpha}=\alpha^{\sim}=\left(P_{v} \overline{\alpha_{v}} P_{v}\right)_{v \in \mathbf{a}}$. Then we can easily verify that $\widetilde{\alpha} \in \prod_{v \in \mathbf{a}} U\left(\varphi_{v}^{\prime}\right), \widetilde{\alpha w}=\widetilde{\alpha} \widetilde{w}, \lambda_{v}(\widetilde{\alpha}, \widetilde{w})=R \overline{\lambda_{v}(\alpha, w)} R_{v}$, and $\mu_{v}(\widetilde{\alpha}, \widetilde{w})=\overline{\mu_{v}(\alpha, w)}$ for every $v \in \mathbf{a}$.

Now we define $J: \mathfrak{Z}^{\psi} \times \mathfrak{Z}^{\varphi} \rightarrow \mathcal{H}_{q+n}^{\mathbf{a}}$ and $I: G_{\mathbf{a}}^{\psi} \times G_{\mathbf{a}}^{\varphi} \rightarrow G_{\mathbf{a}}^{\eta}$ by $J(z, w)=\iota(z, \widetilde{w})$ and $I(\beta, \gamma)=\left[\tau_{v} \beta \tau_{v}^{-1},\left(\sigma_{v} \gamma_{v} \sigma_{v}^{-1}\right)^{\sim}\right]_{R}$ with the symbols of (26.2) and [S97, (6.10.2), (6.10.5), (12.1.4), (22.2.1)].

For $f \in C^{\infty}(\mathcal{H})$ define $\tilde{f}=f^{\sim} \in C^{\infty}(\mathcal{H})$ by $\tilde{f}(w)=\overline{f(\widetilde{w})}$. Then $\left(f \|_{\kappa} \alpha\right)^{\sim}=$ $\widetilde{f} \|_{\kappa} \widetilde{\alpha}$ for every $\alpha \in G_{\mathbf{a}}^{\varphi}$ and $\kappa \in \mathbf{Z}^{\mathbf{a}}$. Define $L_{v}$ and $\mathcal{L}_{v}$ as in the proof of Lemma A8.7 for a fixed $\kappa$, and define a differential operator $\widetilde{L}_{v}$ on $\mathcal{H}$ by $\widetilde{L}_{v} f=\left(L_{v} \widetilde{f}\right)^{\sim}$. Then we easily see that $\widetilde{L}_{v}\left(f \|_{\kappa} \alpha\right)=\left(\widetilde{L}_{v} f\right) \|_{\kappa} \alpha$ for every $\alpha \in G_{\mathbf{a}}^{\varphi}$. Now all differential operators $D$ on $\mathcal{H}_{v}$ such that $D\left(f \|_{\kappa_{v}} \alpha\right)=(D f) \|_{\kappa_{v}} \alpha$ for every $\alpha \in G_{v}^{\varphi}$ form a polynomial ring generated by $r_{v}$ elements $D_{i}, 1 \leq i \leq r_{v}$, such that $D_{i}$ is of degree $2 i$ and $D_{1}=L_{v}$. (See [S90, Theorem 3.6 (3), (4)].) Since $\widetilde{L}_{v}$ is of degree 2, we have $\widetilde{L}_{v}=a L_{v}+b$ with constants $a$ and $b$. Take a nonzero holomorphic function $f$ on $\mathcal{H}_{v}$. Then $\tilde{f}$ is holomorphic, and so $L_{v} f=L_{v} \widetilde{f}=0$, so that $\widetilde{L}_{v} f=0$. Therefore $b$ must be 0 . Thus $\widetilde{L}_{v}=a L_{v}$, and clearly $a \neq 0$. This shows that $L_{v} g=0$ if and only if $\widetilde{L}_{v} g=0$, that is, if and only if $L_{v} \widetilde{g}=0$.

Let $R_{1}=R \| U^{-1}, \rho_{0}(a, b)=\operatorname{det}(b)^{\nu \mathbf{a}}$, and $\omega(a, b)=\operatorname{det}(b)^{m}$. By the same argument as in §A8.12, we find an element $\varepsilon$ of $\mathfrak{U}$ such that $\left(\varepsilon R_{1}^{\rho_{0}}\right)^{\circ}=\left(S_{0}^{\omega}\right)^{\circ}$. By Theorem 14.12 (4), $\pi^{-c} \Delta_{\nu \mathbf{a}}^{p} R$ for some $c \in \mathbf{Z}$ is a $\overline{\mathbf{Q}}$-rational nearly holomorphic function. Define $L_{v}$ on $\mathfrak{Z}^{\varphi}$ as in the present lemma and $\mathcal{L}_{v}$ on $G_{\mathbf{a}}^{\varphi}$ as in the proof of Lemma A8.7. Let $M(z, w)=\pi^{-c} S_{0}(\iota(z, w))$. Since $S_{0}=B\left[\left(\Delta_{\nu \mathbf{a}}^{p} R\right) \| U^{-1}\right]$, formula (29.6) is valid with $M$ in place of $A_{q}^{k} f$. By (12.32) and Theorem 14.9 (2), $\mathcal{N}_{h}^{e^{\prime}}(\Gamma, \overline{\mathbf{Q}})$ is stable under the $L_{v}$ for any congruence subgroup $\Gamma$ of $G^{\varphi}$. Thus $\overline{M(z, w)}$ belongs to a finite-dimensional vector space over $\overline{\mathbf{Q}}$ that is stable under the $L_{v}$. Put $S_{1}(z, w)=S(z, \widetilde{w})$. Then $S_{1}(z, w)=S_{0}(J(z, w))$ and $\overline{M(z, w)}=$
$\pi^{-c} \overline{S_{1}(z, \widetilde{w})}$ Let $\mathcal{X}=\sum_{v \in \mathbf{a}} \sum_{i=0}^{\infty} \mathbf{C} L_{v}^{i} \bar{M}, \mathcal{X}(\overline{\mathbf{Q}})=\sum_{v \in \mathbf{a}} \sum_{i=0}^{\infty} \overline{\mathbf{Q}} L_{v}^{i} \bar{M}$, and $\mathcal{X}^{\prime}=$ $\sum_{v \in \mathbf{a}} \sum_{i=0}^{\infty} \mathbf{C} L_{v}^{i} S_{1}$. Then $\mathcal{X}(\overline{\mathbf{Q}})$ is finite-dimensional over $\overline{\mathbf{Q}}$. For $f \in C^{\infty}\left(\mathfrak{Z}^{\psi} \times \mathfrak{Z}^{\varphi}\right)$ define $f^{\omega} \in C^{\infty}\left(G_{\mathbf{a}}^{\psi} \times G_{\mathbf{a}}^{\varphi}\right)$ by $f^{\omega}(x, y)=\omega\left(M_{(x, y)}(\mathbf{o})\right)^{-1} f(x(\mathbf{o}), y(\mathbf{o}))$. Then $S_{1}^{\omega}=$ $\left(S_{0}^{\omega}\right)^{\circ}=\left(\varepsilon R_{1}^{\rho_{0}}\right)^{\circ}$, which is contained in the unitarizable space by Lemma A8.11. Therefore the $\mathcal{L}_{v}$ are diagonalizable on $\left(\mathcal{X}^{\prime}\right)^{\omega}$, so that the $L_{v}$ are diagonalizable on $\mathcal{X}^{\prime}$. Now $f \mapsto \widetilde{f}$ maps $\mathcal{X}^{\prime}$ (anti-C-linearly) onto $\mathcal{X}$, and hence the $L_{v}$ and $\widetilde{L}_{v}$ are diagonalizable on $\mathcal{X}$. Therefore we have a decomposition $\mathcal{X}(\overline{\mathbf{Q}})=\mathcal{Y}_{0} \oplus \mathcal{Y}$ with $\mathcal{Y}_{0}=\left\{h \in \mathcal{X}(\overline{\mathbf{Q}}) \mid L_{v} h=0\right.$ for every $\left.v \in \mathbf{a}\right\}$ and a vector space $\mathcal{Y}$ spanned by the eigenfunctions $g$ of the $L_{v}$ such that $L_{v} g \neq 0$ for some $v$. Then $\langle g(z, w), f(w)\rangle=$ 0 for every $f \in \mathcal{S}_{h}^{\varphi}$ and every $g \in \mathcal{Y}$ for the same reason as in the proof of Lemma A8.7. Let $T_{1}(z, w)$ be the projection of $\overline{M(z, w)}$ to $\mathcal{Y}_{0}$. Then $\langle\overline{M(z, w)}-$ $\left.T_{1}(z, w), f(w)\right\rangle=0$ for every $f \in \mathcal{S}_{h}^{\varphi}$ and $L_{v} T_{1}=0$ for every $v$. Thus $\widetilde{L}_{v} T_{1}=0$, and so $L_{v} \widetilde{T}_{1}=\left(\widetilde{L}_{v} T_{1}\right)^{\sim}=0$. Since $\widetilde{T}_{1} \in \mathcal{X}^{\prime}$, (A8.12) shows that $\widetilde{T}_{1}$ is holomorphic in $w$, so that $T_{1}$ is holomorphic in $w$. Thus $\pi^{c} T_{1}$ gives the desired $T$ of our lemma.

This page intentionally left blank

## REFERENCES

[F] P. Feit, Poles and residues of Eisenstein series for symplectic and unitary groups, Memoirs, Amer. Math. Soc. 61, No. 346 (1986).
[H] G. Harder, A Gauss-Bonnet formula for discrete arithmetically defined groups, Ann. Sci. École Norm. Sup. $4^{e}$ série, 4, (1971), 409-455.
[K] H. Klingen, Über den arithmetischen Character der Fourierkoeffizienten von Modulformen, Math. Ann. 147 (1962), 176-188.
[M] T. Miyake, On $\mathbf{Q}_{\mathrm{ab}}$-rationality of Eisenstein series of weight 3/2, J. Math. Soc. Japan, 41 (1989), 473-492.
[P] T-y. Pei, Eisenstein series of weight 3/2; I, II, Trans. Amer. Math. Soc. 274 (1982), 573-606, 283 (1984), 589-603.
[R] D. Rohrlich, Nonvanishing of $L$-functions for $G L(2)$, Inv. math. 97 (1989), 381-403.
[S59] G. Shimura, On the theory of automorphic function, Ann. of Math. 70 (1959), 101-144.
[S63] G. Shimura, On analytic families of polarized abelian varieties and automorphic functions, Ann. of Math. 78 (1963), 149-192.
[S64] G. Shimura, On the field of definition for a field of automorphic functions, Ann. of Math. 80 (1964), 160-189.
[S65] G. Shimura, On the field of definition for a field of automorphic functions: II, Ann. of Math. 81 (1965), 124-165.
[S66a] G. Shimura, Moduli and fibre systems of abelian varieties, Ann. of Math. 83 (1966), 294-338.
[S66b] G. Shimura, On the field of definition for a field of automorphic functions: III, Ann. of Math. 83 (1966), 377-385.
[S67] G. Shimura, Construction of class fields and zeta functions of algebraic curves, Ann. of Math. 85, 58-159 (1967).
[S70] G. Shimura, On canonical models of arithmetic quotients of bounded symmetric domains, Ann. of Math. 91 (1970), 144-222; II, 92 (1970), 528-549.
[S71] G. Shimura, Introduction to the arithmetic theory of automorphic functions, Publ. Math. Soc. Japan, No.11, Iwanami Shoten and Princeton Univ. Press, 1971.
[S75] G. Shimura, On some arithmetic properties of modular forms of one and several variables, Ann. of Math. 102 (1975), 491-515.
[S76] G. Shimura, The special values of the zeta functions associated with cusp forms, Comm. pure appl. Math. 29 (1976), 783-804.
[S78a] G. Shimura, The arithmetic of automorphic forms with respect to a unitary group, Ann. of Math. 197 (1978), 569-605.
[S78b] G. Shimura, On certain reciprocity-laws for theta functions and modular forms, Acta math. 141 (1978), 35-71.
[S78c] G. Shimura, On some problems of arithmeticity, Proc. Int. Congress of Math. Helsinki, 1978, 373-379.
[S79] G. Shimura, Automorphic forms and the periods of abelian varieties, J. Math. Soc. Japan, 31 (1979), 561-592.
[S80] G. Shimura, The arithmetic of certain zeta functions and automorphic forms on orthogonal groups, Ann. of Math. 111 (1980), 313-375.
[S81a] G. Shimura, The critical values of certain zeta functions associated with modular forms of half-integral weight, J. Math. Soc. Japan 33 (1981), 649-672.
[S81b] G. Shimura, Arithmetic of differential operators on symmetric domains, Duke M. J. 48 (1981), 813-843.
[S82] G. Shimura, Confluent hypergeometric functions on tube domains, Math. Ann. 260 (1982), 269-302.
[S83] G. Shimura, On Eisenstein series, Duke Math. J. 50 (1983), 417-476.
[S84a] G. Shimura, Differential operators and the singular values of Eisenstein series, Duke Math. J. 51 (1984), 261-329.
[S84b] G. Shimura, On differential operators attached to certain representations of classical groups, Inv. math. 77 (1984), 463-488.
[S85a] G. Shimura, On Eisenstein series of half-integral weight, Duke Math. J. 52 (1985), 281-324.
[S85b] G. Shimura, On the Eisenstein series of Hilbert modular groups, Revista Mat. Iberoamer. 1, No. 3 (1985), 1-42.
[S86] G. Shimura, On a class of nearly holomorphic automorphic forms, Ann. of Math. 123 (1986), 347-406.
[S87a] G. Shimura, Nearly holomorphic functions on hermitian symmetric spaces, Math. Ann. 278 (1987), 1-28.
[S87b] G. Shimura, On Hilbert modular forms of half-integral weight, Duke M. J. 55 (1987), 765-838.
[S88] G. Shimura, On the critical values of certain Dirichlet series and the periods of automorphic forms, Inv. math. 93 (1988), 1-61.
[S90] G. Shimura, Invariant differential operators on hermitian symmetric spaces, Ann. of Math. 132 (1990), 237-272.
[S91] G. Shimura, The critical values of certain Dirichlet series attached to Hilbert modular forms, Duke Math. J. 63 (1991), 557-613.
[S93] G. Shimura, On the transformation formulas of theta series, Amer. J. of Math. 115 (1993), 1011-1052.
[S94a] G. Shimura, Euler products and Fourier coefficients of automorphic forms on symplectic groups, Inv. math. 116 (1994), 531-576.
[S94b] G. Shimura, Differential operators, holomorphic projection, and singular forms, Duke Math. J. 76 (1994), 141-173.
[S95a] G. Shimura, Eisenstein series and zeta functions on symplectic groups, Inv. math. 119 (1995), 539-584.
[S95b] G. Shimura, Zeta functions and Eisenstein series on metaplectic groups, Inv. math. 121 (1995), 21-60.
[S96] G. Shimura, Convergence of zeta functions on symplectic and metaplectic groups, Duke Math. J. 82 (1996), 327-347.
[S97] G. Shimura, Euler Products and Eisenstein series, CBMS Regional Conference Series in Mathematics, No.93, Amer. Math. Soc., 1997.
[S98] G. Shimura, Abelian varieties with complex multiplication and modular functions, Princeton University Press, 1998.
[S99] G. Shimura, The number of representations of an integer by a quadratic form, Duke Math. J. 100 (1999), 59-92.
[Si] C. L. Siegel, Gesammelte Abhandlungen, I-III, 1966; IV, 1979, Springer.
[St] J. Sturm, The critical values of zeta functions associated to the symplectic group, Duke Math. J. 48 (1981), 327-350.
[W46] A. Weil, Foundations of algebraic geometry, Amer. Math. Soc. 1946, 2nd ed. 1962.
[W48] A. Weil, Variétés abéliennes et courbes algébriques, Hermann, Paris, 1948.
[W56] A. Weil, The field of definition of a variety, Amer. J. Math., 78 (1956), 509-524.
[W58] A. Weil, Introduction à l'étude des variétés Kählériennes, Hermann, Paris, 1958.
[W64] A. Weil, Sur certain groupes d'opérateurs unitaires, Acta Math. 111 (1964), 143-211.

This page intentionally left blank

## Index

a, 9
$\mathcal{A}_{k}, 32$
$\mathcal{A}_{\omega}, 32$
b, 30
$\operatorname{det}_{i}(x), 91$
$\mathbf{e}_{\mathbf{a}}, 10,32$
$\mathbf{e}_{\mathbf{A}}, \mathbf{e}_{\mathbf{h}}, \mathbf{e}_{p}, \mathbf{e}_{v}, 10$
$\mathbf{e}_{\mathbf{a}}^{n}, 32$
$\mathbf{e}_{\mathbf{A}}^{n}, \mathbf{e}_{\mathbf{h}}^{n}, \mathbf{e}_{v}^{n}, 127$
$\tilde{G}, 20$
$G_{0}, 21$
$\mathcal{G}_{+}, 54$
g-lattice, 9
$\Gamma_{n}^{\ell}, 135$
h, 9
$\mathcal{H}_{\mathrm{CM}}, 76$
$I_{K}, 77$
$\mathrm{il}_{3}(x), 12$
$j_{\alpha}(z), 19,30$
$j^{k}, 130$
$J_{K}, 77$
[k], 130
$k_{\Omega}, 58$
$\mathfrak{K}$ (field), 47
$\mathfrak{K}$ (group), 110
$\mathfrak{K}_{0}, 89$
$M(\alpha, z), 26,30$
$M l_{p}(),$,
$\mathcal{M}_{k}, 31$
$\mathcal{M}_{k}^{\varphi}, 208$
$\mathcal{M}_{\omega}, 31$
$\mathfrak{M}, 129,254$
$\nu(x), 11$
$\nu_{0}(x), 10$
$\nu_{\mathrm{b}}(\xi), 161$
$p_{Y}(\alpha, \beta), 78$
$\mathfrak{p}(w), 78$
$\mathfrak{P}_{*}(w), 82$
$S_{p}(),$,
$\mathcal{S}_{k}, 35,208$
$\mathcal{S}_{k}^{\varphi}, 208$
$\mathcal{S}_{\omega}, 34$
$\mathfrak{T}(s, \mathbf{f}, \chi), 173$
$\mathfrak{U}, 285$
v, 9
$\mathbf{Z}_{\mathbf{h}}^{\times}, 50$
$\mathcal{Z}(s, \mathbf{f}, \chi), 173,176$
$3^{\varphi}, 208$
abelian variety, 13
adelization, 9
anisotropic, 7
automorphic form, 32
automorphic function, 32
canonical model, 58
Case SP, 17
Case UB, 17
Case UT, 17
CM-algebra, 17
CM-field, 12
CM-type, 12
component of an automorphic form, 166
congruence subgroup, 27
cusp condition, 33, 103
cusp form, 34,73
denominator ideal, 10
factor of automorphy, 18,19
factor of automorphy of weight $1 / 2,40,254$
field of moduli, 16
Gauss sum, 156
generic, 46
half-integral weight, 41
Hecke character, 129
integral weight, 41
involution of the endomorphism algebra, 15
isotropic, 7
lattice in a real vector space, 13
Möbius function, 163
metaplectic group, 252
modular form, 32
modular form of half-integral weight, 41
nearly holomorphic function, 100
PEL-type, 23
period symbol, 77
polarization, 14
polarized abelian variety, 14
primitive, 10
quasi-representation, 72
rationality of automorphic forms, 34, 82
reduced expression, 10
reflex field, 12, 60
reflex of a CM-type, 60
Riemann form, 13
smallest field of definition, 15
unitarizable $\mathfrak{U}$-module, 285

Written by one of the leading experts, venerable grandmasters, and most active contributors ... in the arithmetic theory of automorphic forms ... the new material included here is mainly the outcome of his extensive work ... over the last eight years . . a very careful, detailed introduction to the subject ... this monograph is an important, comprehensively written and profound treatise on some recent achievements in the theory:
-Zentralblatt MATH
The main objects of study in this book are Eisenstein series and zeta functions associated with Hecke eigenforms on symplectic and unitary groups. After preliminaries-including a section, "Notation and Terminology"-the first part of the book deals with automorphic forms on such groups. In particular, their rationality over a number field is defined and discussed in connection with the group action; also the reciprocity law for the values of automorphic functions at CM-points is proved. Next, certain differential operators that raise the weight are investigated in higher dimension. The notion of nearly holomorphic functions is introduced, and their arithmeticity is defined. As applications of these, the arithmeticity of the critical values of zeta functions and Eisenstein series is proved.
Though the arithmeticity is given as the ultimate main result, the book discusses many basic problems that arise in number-theoretical investigations of automorphic forms but that cannot be found in expository forms. Examples of this include the space of automorphic forms spanned by cusp forms and certain Eisenstein series, transformation formulas of theta series, estimate of the Fourier coefficients of modular forms, and modular forms of halfintegral weight. All these are treated in higher-dimensional cases. The volume concludes with an Appendix and an Index.

The book will be of interest to graduate students and researchers in the field of zeta functions and modular forms.

