

A Poisson formula and exponential sums¹⁾

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Introduction. We shall first recall our Poisson formula in [4] with precise definitions: let k denote a global field; we shall denote by the subscripts A and k the adelization relative to k and the taking of k -rational points, respectively. We shall denote by k_A and k_A^\times the adèle and idele groups of k ; we shall fix a non-trivial character ψ of k_A/k . Let X denote an affine n -space over k with x_1, x_2, \dots, x_n as its affine coordinates; let dx denote the exterior product of their differentials and $|dx|_A$ the Haar measure on X_A such that X_A/X_k is of measure 1. We shall fix a form $f(x)$ of degree $m \geq 2$ in x_1, x_2, \dots, x_n with coefficients in k ; we shall assume that $\text{char}(k)$ does not divide m , $f(x)$ is "strongly non-degenerate", i.e., the projective hypersurface defined by $f(x)=0$ is non-singular, and $n > 2m$. Let $U(i)$ denote the set of simple points of the hypersurface $f(x)=i$ and $\theta_i(x)$ the residue of $(f(x)-i)^{-1}dx$ along $U(i)$; then if i is in k , the gauge-form θ_i on $U(i)$ gives rise to a measure $|\theta_i|_A$ on $U(i)_A$. Finally let $\mathcal{S}(X_A)$ denote the Schwartz-Bruhat space of X_A and take Φ from $\mathcal{S}(X_A)$. Then we have the following identity:

$$\sum_{i \in k} \int_{U(i)_A} \Phi |\theta_i|_A = \sum_{i^* \in k} \int_{X_A} \psi(i^* f(x)) \Phi(x) |dx|_A;$$

both sides have dominant series if Φ is restricted to a compact subset of $\mathcal{S}(X_A)$. We shall add the Dirac measure δ_0 to the so-defined tempered measure on X_A and denote the new measure by E ; then $E(\Phi) - \Phi(0)$ is given by either side of the above identity.

On the other hand

$$E'(\Phi) = \sum_{\xi \in X_k} \Phi(\xi)$$

defines a tempered measure E' on X_A independently of $f(x)$ for every n . We observe that $\text{Supp}(E')$, the support of E' , is the discrete subset X_k of X_A while $\text{Supp}(E)$ is in general the union of $f^{-1}(i)_A$ for all i in k ; $\text{Supp}(E)$ becomes the union of $\{0\}$ and $f^{-1}(i)_A = U(i)_A$ for all i in k^\times if $U(0)_A$ is empty. The *main objective* of this paper is to prove the following remarkable relation between E and E' :

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We convert $G_A = k_A \times k_A^\times$ into a (locally compact) group as

$$(u, t) \cdot (u', t') = (u + t^m u', tt')$$

and for every $g = (u, t)$ in G_A we put

$$(U(g)\Phi)(x) = |t|_A^{(1/2)n} \phi(uf(x))\Phi(tx),$$

in which $|t|_A$ is the modulus of t . Then $(E' - E)(U(g)\Phi)$ becomes a G_k -invariant continuous function on G_A and it vanishes at every k -rational boundary point of G_A to the order $s-1$ where $s = n/2m$. This is a generalization of a classical theorem stating that the theta series defined by a quadratic form behaves up to lower order terms like the corresponding Eisenstein series at every cusp. We recall that if $f(x)$ is a quadratic form, then the above G_A can be embedded in the "metaplectic group" which is at most a two-sheeted covering group of $(SL_2)_A$ and which contains $(SL_2)_k$ as a subgroup and that $E' - E$ is $(SL_2)_k$ -invariant; in view of the reduction theory for $(SL_2)_A$ the above-recalled theorem then implies the boundedness of $(E' - E)(U(g)\Phi)$ on G_A . We refer to Weil [12], Chap. V for the details.

The theory of metaplectic groups for higher degree forms is not yet available. However our result seems to indicate that a boundedness as above holds in some form. We have included a remark that a weak form of the boundedness implies a generalization of the Hasse-Minkowski theorem. We have also included as an appendix an observation showing that a naive generalization of the metaplectic group might give something different from SL_2 . It is hoped that this paper serves as a guide to the future arithmetic theory of higher degree forms.

1. Reduction theory for G_A . We shall denote by p a (not necessarily non-archimedean) valuation on the global field k , by k_p the corresponding completion of k , and by $|\cdot|_p$ the usual absolute value on k_p . If t_p is an element of k_p^\times , then $|t_p|_p$ is the modulus of t_p , i.e., the rate of measure change in k_p under the multiplication by t_p . Moreover the modulus $|t|_A$ of an element $t = (t_p)_p$ of k_A^\times becomes the product of all $|t_p|_p$. We recall that $|t|_A = 1$ if t is in k^\times ; we shall use this "product-formula" all the time.

Since $G_A = k_A \times k_A^\times$, every element g of G_A is a pair (u, t) of an element u of k_A and an element t of k_A^\times ; g is in G_k if u is in k and t in k^\times . We put

$$\mathbf{t}(u) = U(u, 1), \quad \mathbf{d}(t) = U(0, t);$$

then we get

$$(\mathbf{t}(u)\Phi)(x) = \phi(uf(x))\Phi(x), \quad (\mathbf{d}(t)\Phi)(x) = |t|_A^{(1/2)n} \Phi(tx)$$

for every Φ in $\mathcal{S}(X_A)$ and

$$U(u, t) = \mathbf{t}(u)\mathbf{d}(t) = \mathbf{d}(t)\mathbf{t}(t^{-m}u).$$

We recall that the correspondence $(g, \Phi) \rightarrow U(g)\Phi$ does define a mapping $G_A \times \mathcal{S}(X_A) \rightarrow \mathcal{S}(X_A)$ and that it is continuous; such a continuity problem was examined quite generally by Levin [7]. Moreover $g \rightarrow U(g)$ uniquely extends to a unitary representation of G_A in $L^2(X_A)$. Since E, E' are tempered distributions on X_A , both $E(U(g)\Phi)$ and $E'(U(g)\Phi)$ become complex-valued continuous functions on $G_A \times \mathcal{S}(X_A)$. We observe that E, E' are G_k -invariant, i.e., $E(U(g)\Phi) = E(\Phi)$, $E'(U(g)\Phi) = E'(\Phi)$ for every g in G_k and Φ in $\mathcal{S}(X_A)$; the verification is straightforward.

We shall construct a convenient subset of G_A whose left translates by elements of G_k cover the whole G_A ; such a procedure is called a reduction theory for G_A . We shall first recall the reduction theories for k_A and k_A^\times : if k is a number field, we shall denote by S_∞ the set of all archimedean valuations on k ; then k_p is an \mathbf{R} -field, i.e., $k_p = \mathbf{R}$ or \mathbf{C} , for every p in S_∞ . We shall denote by \mathfrak{o} the ring of integers of k , i.e., the integral closure of \mathbf{Z} in k . If k is a function field, we shall denote by \mathbf{F}_{q_0} the algebraic closure in k of the prime field; we then choose a valuation p_∞ such that the corresponding residue class field \mathbf{F}_q is of degree at least twice the genus of k over \mathbf{F}_{q_0} , and we take $\{p_\infty\}$ as S_∞ . According to the Riemann-Roch theorem there exists an element T of k with p_∞ as its polar divisor, and we have

$$k_{p_\infty} = \mathbf{F}_q((T^{-1})).$$

This time we shall denote by \mathfrak{o} the integral closure of $\mathbf{F}_{q_0}[T]$ in k .

In either case if k_p is a p -field, i.e., if k_p is not an \mathbf{R} -field, we shall denote by X_p^0 the compact open subgroup $(k_p^0)^n$ of $X_p = k_p^n$, in which k_p^0 is the maximal compact subring of k_p , and by X_0 the restricted product of X_p relative to X_p^0 for all p not in S_∞ ; then

$$X_0^0 = \prod_{p \notin S_\infty} X_p^0$$

becomes a compact open subgroup of X_0 . We shall denote by X_∞ the product of X_p for all p in S_∞ ; then we have

$$X_A = X_0 \times X_\infty;$$

accordingly we shall express an element x of X_A as (x_0, x_∞) . We shall denote by X_0 the subgroup of X_k consisting of \mathfrak{o} -rational points of X . If we diagonally embed X_k in X_A , then we have

$$X_A = X_k + (X_0^0 \times X_\infty), \quad X_0 = X_k \cap (X_0^0 \times X_\infty);$$

moreover the projection of X_0 to X_∞ becomes a lattice in X_∞ ; we shall denote X_0, X_0^0, X_∞ for $n=1$ by k_0, k_0^0, k_∞ . This is the reduction theory for X_A and in particular for k_A .

If k is a number field and τ is in \mathbf{R}_+^\times , the multiplicative group of positive real numbers, then we shall denote by a_τ the element of k_A^\times defined by $a_\tau = (1, \tau)$, i.e., by the condition that non-archimedean components are 1 and archimedean components are τ . If k is a function field, we replace \mathbf{R}_+^\times by the cyclic subgroup of $(k_{p_\infty})^\times$ generated by T and define a_τ similarly as above. Then in either case we have

$$k_A^\times = k^\times \cdot \{a_\tau\}_\tau \cdot \text{compact};$$

this is the reduction theory for k_A^\times and it implies the finiteness of the ideal class group and the theorem of units of \mathfrak{o} . All these are well known as the Iwasawa-Tate theory; we refer to [11], [13] for the details.

The reduction theory for G_A follows from what we have recalled; we have

$$G_A = G_k \cdot \{(u, a_\tau); u_0 = 0, u_\infty \in \text{compact}\} \cdot \text{compact}.$$

The proof is as follows: let (u, t) denote an arbitrary element of G_A ; then we can write

$$t = i a_\tau c, \quad i^{-m} u = i^{-m} i^* + (v_0, c^*)$$

with i in k^\times , c in a fixed compact subset of k_A^\times , i^* in k , v_0 in k_0^0 , and c^* in a fixed compact subset of k_∞ ; and we get

$$(u, t) = (i^*, i)((0, c^*), a_\tau)((v_0, 0), c).$$

This proves the assertion.

We observe that the middle space can be injectively mapped to $k_\infty \times \mathbf{R}_+^\times$. We take the one point compactification of $k_\infty \times (\mathbf{R}_+^\times \cup \{0\})$ and regard the boundary of $k_\infty \times \mathbf{R}_+^\times$ as the boundary of G_A . In particular points of $k \times \{0\}$ and the point at infinity will be considered as k -rational boundary points of G_A . We shall examine the behavior of $(E' - E)(U(g)\Phi)$ as g approaches such boundary points; for the exact statement of our results we refer to § 4, Th. 3 and Th. 4.

2. A dominant series. If $f(x)$ is a quadratic form and k is a number field, the series

$$E(\Phi) = \Phi(0) + \sum_{i^* \in k} \int_{X_A} \phi(i^* f(x)) \Phi(x) |dx|_A$$

appeared (as a special case) in Weil [12], Chap. IV; he called it the Eisenstein-Siegel series. We recall that if $k=\mathbf{Q}$ and $f(x_\infty)$ is a positive-definite quadratic form, then $E(U(g)\Phi)$ becomes the classical holomorphic Eisenstein series for a suitable Φ . We shall show in the general case where $f(x)$ is a higher degree form that $E(U(g)\Phi)$ has a certain Eisenstein series as a dominant series; we shall first construct an equivalent dominant series.

If $g=(u, t)$ is in G_A , we get

$$\begin{aligned} & \int_{X_A} \phi(i^*f(x))(U(g)\Phi)(x)|dx|_A \\ &= |t|_A^{(1/2)n-n} \cdot \int_{X_A} \phi((i^*+u)t^{-m}f(x))\Phi(x)|dx|_A, \end{aligned}$$

hence its absolute value is at most equal to

$$c(\Phi)|t|_A^{(1/2)n} \cdot \prod_p \max(|u_p+i^*|_p, |t_p|_p^m)^{-n/m},$$

in which $c(\Phi)$ is independent of i^* , t , and u ; cf. [4]. Actually if we restrict Φ to a compact subset of $\mathcal{S}(X_A)$, then we can replace $c(\Phi)$ by a constant independent of Φ ; cf. 2 bis. We shall show that $E(U(g)\Phi)$ has

$$\text{const. } |t|_A^{(1/2)n} (1 + \sum_{i^* \in k} \prod_p \max(|u_p+i^*|_p, |t_p|_p^m)^{-n/m})$$

as a dominant series; we have only to show that this series is convergent.

For the sake of simplicity we shall write $|u|_p$, etc. instead of $|u_p|_p$, etc. We shall also use the following notation: if ϕ and ϕ' are complex-valued functions on a set X such that

$$|\phi(x)| \leq c \cdot \phi'(x)$$

for every x in X , in which c is a constant, we write $\phi(x) < \phi'(x)$ or $\phi'(x) > \phi(x)$ on X ; if we have both $\phi(x) < \phi'(x)$ and $\phi(x) > \phi'(x)$ on X , we write $\phi(x) \succ < \phi'(x)$ on X . If there is no ambiguity about the set X , we sometimes omit it. Finally we shall use this notation to define a family of subsets of a given set. The following lemma will settle the convergence problem:

LEMMA 1. *Suppose that $s > 1$ and (u, v) is in $k_A \times k_A^\times$; then we have*

$$\sum_{i^* \in k} \prod_p \max(|u+i^*|_p, |v|_p)^{-2s} < |v|_A^{1-2s}$$

on the subset of $k_A \times k_A^\times$ defined by $|v|_A > 1$.

PROOF. By the reduction theory for k_A^\times we can write v as $ia_\tau c$, in which i is in k^\times and c is in a fixed compact subset of k_A^\times ; then we get $|v|_A \succ < |a_\tau|_A$ and

$$\prod_p \max(|u+i^*|_p, |i a_\tau c|_p) > < \prod_p \max(|(i^{-1}u)+(i^{-1}i^*)|_p, |a_\tau|_p).$$

Therefore for our purpose we may assume that $v=a_\tau$ and $|v|_A > 1$; this implies $|v|_p > 1$ for every p in S_∞ and $|v|_p = 1$ for every other p . By the reduction theory for k_A we may assume that $|u|_p < 1$ for every p in S_∞ and $|u|_p \leq 1$ for every other p . Then we get

$$\max(|u+i^*|_p, |v|_p) > < \max(|i^*|_p, |v|_p)$$

for every p in S_∞ and

$$\max(|u+i^*|_p, |v|_p) = \max(|i^*|_p, |v|_p)$$

for every other p . Therefore we get

$$\begin{aligned} \sum_{i^* \in k} \prod_p \max(|u+i^*|_p, |v|_p)^{-2s} \\ > < \sum_{i^* \in k} \prod_p \max(|i^*|_p, |v|_p)^{-2s}; \end{aligned}$$

the rest of the proof is as follows:

Since the class number of k or rather of \mathfrak{o} is finite, we can express i^* as ab^{-1} with a, b in \mathfrak{o} , $b \neq 0$ such that

$$\prod_{p \in S_\infty} \max(|a|_p, |b|_p) > < 1.$$

Furthermore, in view of the theorem of units, we may assume that b is among a fixed set, say R , of representatives of orbits in $\mathfrak{o} - \{0\}$ by the group of units of \mathfrak{o} with the property that $|b|_p > 1$ for every p in S_∞ . Then we get

$$\begin{aligned} \sum_{i^* \in k} \prod_p \max(|i^*|_p, |v|_p)^{-2s} \\ < \sum_{b \in R} \sum_{a \in \mathfrak{o}} \prod_{p \in S_\infty} \max(|a|_p, |bv|_p)^{-2s}; \end{aligned}$$

and we have

$$\sum_{a \in \mathfrak{o}} \prod_{p \in S_\infty} \max(|a|_p, |bv|_p)^{-2s} < \prod_{p \in S_\infty} |bv|_p^{1-2s}.$$

This is the crucial point of our proof: if k is a number field, it follows from Lemma 12 in [3]; and if k is a function field, it follows from a counterpart of "Lemma 12" proved in [4], p. 227; cf. also Mars [8], pp. 133-136. At any rate, since $|v|_p = 1$ for every p not in S_∞ , we get

$$\prod_{p \in S_\infty} |v|_p = |v|_A.$$

Moreover we have

$$\sum_{b \in R} \left(\prod_{p \in S_\infty} |b|_p \right)^{1-2s} < \infty$$

because this series is a partial sum of the Dedekind zeta series of \mathfrak{o} evaluated at $2s-1 > 1$. q.e.d.

2 bis. We shall outline a method to obtain a dominant series for any series such as $E(\Phi)$; we shall first make some of our notations precise: k is a global field, ϕ is a non-trivial character of k_A/k , and ϕ_p is the product of the canonical injection $k_p \rightarrow k_A$ and ϕ ; X is an affine n -space over k , $[x, y]$ is a symmetric non-degenerate bilinear form on $X \times X$ defined over k , and $|dx|_p$ is the Haar measure on X_p which is autodual relative to the bicharacter $\phi_p([x, y])$ of $X_p \times X_p$; the restricted product measure $|dx|_A$ of all $|dx|_p$ then becomes the Haar measure on X_A such that X_A/X_k is of measure 1.

Let $f(x)$ denote, for a moment, an arbitrary polynomial in the affine coordinates x_1, x_2, \dots, x_n of X with coefficients in k and suppose that

$$(*) \quad \left| \int_{X_p} \phi_p(i^* f(x)) \Phi(x) |dx|_p \right| \leq c(\Phi) \cdot \max(|i^*|_p, 1)^{-\sigma}$$

for every i^* in k_p and Φ in $\mathcal{S}(X_p)$, the Schwartz-Bruhat space of X_p ; it is understood that $c(\Phi)$ is independent of i^* and σ is independent of i^* and Φ . Then for any compact subset C of $\mathcal{S}(X_p)$ we can replace $c(\Phi)$ by a constant which works for all Φ in C . The proof depends on the following fact: "Let B denote a subset of the space $\mathcal{S}(X_p)'$ of tempered distributions on X_p which is bounded in the sense that

$$\sup_{T \in B} |T(\Phi)| < \infty$$

for each Φ in $\mathcal{S}(X_p)$; then we have

$$\sup_{T \in B, \Phi \in C} |T(\Phi)| < \infty."$$

(If k_p is a p -field, this follows from the fact that C is contained in a finite dimensional subspace of $\mathcal{S}(X_p)$. If k_p is an \mathbf{R} -field, then $\mathcal{S}(X_p)$ is a complete metric space; and we have only to apply the Baire category theorem.) If we put

$$T_{i^*}(\Phi) = \max(|i^*|_p, 1)^\sigma \cdot \int_{X_p} \phi_p(i^* f(x)) \Phi(x) |dx|_p$$

for every Φ in $\mathcal{S}(X_p)$, we get a subset $\{T_{i^*}\}_{i^*}$ of $\mathcal{S}(X_p)'$ parametrized by k_p .

The assumption (*) shows that this subset is bounded; the rest follows from the fact that we have recalled.

We can go one step further: "Let B, B' denote bounded subsets of $\mathcal{S}(X_p)'$, $\mathcal{S}(X_{p'})'$, respectively; then the set of $T \otimes T'$ for all T in B and T' in B' is bounded in $\mathcal{S}(X_p \times X_{p'})'$." (This becomes trivial if either k_p or $k_{p'}$ is a p -field; and if both k_p and $k_{p'}$ are \mathbf{R} -fields, then it can be proved by applying twice the above-recalled fact.) Therefore if we have (*) for every p , we can take a finite product. More precisely suppose that S is a finite set of valuations on k and let X_S, k_S denote the products of X_p, k_p for all p in S ; then we have

$$\left| \int_{X_A} \left(\prod_{p \in S} \phi_p(i_p^* f(x_p)) \right) \Phi(x) \left(\otimes_{p \in S} |dx_p|_p \right) \right| \leq c \cdot \prod_{p \in S} \max(|i_p^*|_p, 1)^{-\sigma_p}$$

for every $i^*=(i_p^*)_p$ in k_S and Φ in a compact subset of $\mathcal{S}(X_S)$.

Finally for any non-archimedean valuation p let k_p^0 denote the maximal compact subring of k_p and Φ_p the characteristic function of $X_p^0=(k_p^0)^n$; suppose that we can choose $c(\Phi_p)$'s so that their product is convergent; then we have

$$\left| \int_{X_A} \phi(i^* f(x)) \Phi(x) |dx|_A \right| \leq c' \cdot \prod_p \max(|i_p^*|_p, 1)^{-\sigma_p}$$

for every i^* in k_A and Φ in a compact subset of $\mathcal{S}(X_A)$. Therefore if the series

$$\sum_{i^* \in k} \prod_p \max(|i_p^*|_p, 1)^{-\sigma_p}$$

is convergent, by adding 1 and multiplying a suitable constant we will get a dominant series for $E(\Phi)$. We have shown in [4] that this is the case with $\sigma_p = n/m$ if $f(x)$ is a form of degree $m \geq 2$ such that $\text{char}(k)$ does not divide m , the projective hypersurface defined by $f(x)=0$ is non-singular, and $n > 2m$.

3. Eisenstein series. We shall introduce an Eisenstein series which will serve as a dominant series for $E(U(g)\Phi)$ and examine some of its properties; we shall first make preliminary observations.

Let K denote a local field; K is either an \mathbf{R} -field or a p -field. We embed K in a division ring $L=K+Kw$ as follows: if K is an \mathbf{R} -field, we have a unique choice for L ; it is \mathbf{C} or \mathbf{H} , the Hamilton quaternion algebra, according as K is \mathbf{R} or \mathbf{C} . We shall normalize w as $w^2=-1$. If K is a p -field, we shall take as L an unramified quadratic extension of K ; it exists and is unique. Let K^0, L^0 denote the maximal compact subrings of K, L , respectively; then we shall take as w an element of L such that $L^0=K^0+K^0w$. In either case we shall denote

by Nz the norm of an element z of L ; N gives rise to a continuous homomorphism of L^\times to K^\times .

Let

$$\sigma = \begin{pmatrix} \alpha & \beta \\ \gamma & \delta \end{pmatrix}, \quad z = u + vw$$

denote arbitrary elements of $GL_2(K)$, $Z = K + K^\times w$, respectively; then we have $\gamma z + \delta \neq 0$. Therefore

$$\sigma \cdot z = (\alpha z + \beta)(\gamma z + \delta)^{-1} = u' + v'w$$

is defined and we get

$$v' = v(\sigma \cdot z) = \det(\sigma)v \cdot N(\gamma z + \delta)^{-1} \neq 0;$$

hence $\sigma \cdot z$ is also in Z . Moreover the mapping $GL_2(K) \times Z \rightarrow Z$ defined by $(\sigma \cdot z) \rightarrow \sigma \cdot z$ is continuous and it gives a group action of $GL_2(K)$ on Z ; the action is transitive because

$$\begin{pmatrix} v & u \\ 0 & 1 \end{pmatrix} \cdot w = u + vw$$

for every $z = u + vw$ in Z .

We shall go back to the global field k and take k_p as K ; and we shall denote the above w by w_p . For the sake of completeness (and clarity) we shall prove the following lemma:

LEMMA 2. Let u_p, v_p denote arbitrary elements of k_p ; then we have

$$|N(u_p + v_p w_p)|_p > \max(|u_p|_p, |v_p|_p)^2$$

if k_p is an \mathbf{R} -field and

$$|N(u_p + v_p w_p)|_p = \max(|u_p|_p, |v_p|_p)^2$$

if k_p is a p -field.

PROOF. If a, b are non-negative real numbers, we certainly have

$$\max(a, b) \leq a + b \leq 2 \cdot \max(a, b);$$

this takes care of the \mathbf{R} -field case. In the p -field case, if a, b are elements of $k_p(w_p)$ satisfying $|a|_p \neq |b|_p$, we have

$$|a + b|_p = \max(|a|_p, |b|_p).$$

Therefore if $|u_p|_p \neq |v_p|_p$, we have the equality in question. We shall assume

that $|u_p|_p = |v_p|_p$; then by homogeneity we may further assume that $|u_p|_p = |v_p|_p = 1$. And we have only to show that $|N(u_p + v_p w_p)|_p = 1$. By our choice of w_p the residue classes of 1, w_p are linearly independent over the residue class field of k_p . Since u_p, v_p are units, the residue class of $u_p + v_p w_p$ is different from 0; hence $u_p + v_p w_p$ is a unit, and hence $|N(u_p + v_p w_p)|_p = 1$. q.e.d.

We choose w_p for each p and put

$$w = (w_p)_p, \quad \mathcal{X} = k_A + k_A^\times w;$$

with the obvious topology \mathcal{X} become a locally compact space. We shall denote an element of \mathcal{X} by $z = u + vw = (u_p + v_p w_p)_p$ and define Nz componentwise, i.e., as $(Nz)_p = Nz_p$ for every p ; then we see by Lemma 2 that N gives a continuous mapping of \mathcal{X} to k_A^\times . Furthermore we have

$$|N(z + i^*)|_A < \max(|u + i^*|_p, |v|_p)^2$$

on $\mathcal{X} \times k$. Therefore if s is any complex number satisfying $\text{Re}(s) > 1$, we see by Lemma 1 that the series

$$E(z, s) = |v|_A^s \left(1 + \sum_{i^* \in k} |N(z + i^*)|_A^{-s} \right)$$

is absolutely convergent; this is the Eisenstein series useful for our purpose. The following theorem follows from what we have observed:

THEOREM 1. *Consider the mapping $G_A \rightarrow \mathcal{X}$ defined by $g = (u, t) \rightarrow z = u + vw$ where $v = t^m$; then we have*

$$\int_{X_A} \phi(i^* f(x)) (U(g)\Phi)(x) dx|_A < (|v|_A |N(z + i^*)|_A^{-1})^{n/2m}$$

on $G_A \times k$ if Φ is restricted to a compact subset of $\mathcal{S}(X_A)$. Therefore a constant multiple of $E(u + t^m w, n/2m)$ serves as a dominant series for $E(U(g)\Phi)$ on G_A if Φ is restricted as above. Moreover for any $s > 1$ we have

$$E(z, s) - |v|_A^s < |v|_A^{1-s}$$

on the subset of \mathcal{X} defined by $|v|_A > 1$.

We shall examine the invariance property of $E(z, s)$ under the action of $GL_2(k) = (GL_2)_k$ on \mathcal{X} ; we shall first define the action of $(GL_2)_A$ on \mathcal{X} : let $\sigma = (\sigma_p)_p, z = u + vw$ denote arbitrary elements of $(GL_2)_A, \mathcal{X}$, respectively; we define $\sigma \cdot z = u' + v'w$ componentwise, i.e., as

$$(\sigma \cdot z)_p = u'_p + v'_p w_p = \sigma_p \cdot z_p$$

for every p . We shall show that $\sigma \cdot z$ is also in \mathcal{X} : we have only to show that

if k_p is a p -field, σ_p is in $GL_2(k_p^0)$, and $|u|_p \leq 1$, $|v|_p = 1$, then $|u'|_p \leq 1$, $|v'|_p = 1$. Let $(\gamma \ \delta)$ denote the second row of σ ; then $(\sigma_p \cdot z_p)N(\gamma_p z_p + \delta_p)$ is certainly in $k_p^0 + k_p^0 w_p$. Moreover, since either γ_p or δ_p is a unit of k_p^0 , we get

$$|N(\gamma_p z_p + \delta_p)|_p = \max(|\gamma u + \delta|_p, |\gamma v|_p)^2 = 1;$$

this implies $|u'|_p \leq 1$ and

$$|v'|_p = |\det(\sigma)|_p |v|_p |N(\gamma z + \delta)|_p^{-1} = 1.$$

We observe that the action of $(GL_2)_A$ on \mathcal{X} thus defined is continuous and transitive.

THEOREM 2. *The Eisenstein series $E(z, s)$ is $GL_2(k)$ -invariant, i.e., $E(\sigma \cdot z, s) = E(z, s)$ for every σ in $GL_2(k)$.*

PROOF. We let GL_2 act on the projective line P_1 as

$$\begin{pmatrix} \alpha & \beta \\ \gamma & \delta \end{pmatrix} \cdot i^* = (\alpha i^* + \beta)(\gamma i^* + \delta)^{-1}$$

with the obvious understanding in the case where $i^* = \infty$; then we get a transitive action of $GL_2(k)$ on $P_1(k) = k \cup \{\infty\}$. Let σ denote an arbitrary element of $GL_2(k)$ with coefficients $\alpha, \beta, \gamma, \delta$ and z an element of \mathcal{X} ; then for every i^* in k such that $-\gamma i^* + \alpha \neq 0$, i.e., $\sigma^{-1} \cdot i^* \neq \infty$, we get

$$\sigma \cdot z - i^* = (-\gamma i^* + \alpha)(z - \sigma^{-1} \cdot i^*)(\gamma z + \delta)^{-1}.$$

Moreover if $-\gamma i^* + \alpha = 0$, then $\gamma \neq 0$ and we get

$$\sigma \cdot z - i^* = -\det(\sigma) \gamma^{-1} (\gamma z + \delta)^{-1}.$$

Therefore we get

$$\begin{aligned} E(\sigma \cdot z, s) &= |v(\sigma \cdot z)|_A^s (1 + \sum_{i^* \in k} |N(\sigma \cdot z - i^*)|_A^{-s}) \\ &= |v|_A^s (|N(\gamma z + \delta)|_A^{-s} + \sum_{i^* \neq \infty, \sigma \cdot \infty} |N(z - \sigma^{-1} \cdot i^*)|_A^{-s} + 1) \\ &= |v|_A^s (1 + \sum_{j^* \in k} |N(z - j^*)|_A^{-s}) \end{aligned}$$

if $\sigma \cdot \infty \neq \infty$; and similarly $E(\sigma \cdot z, s) = E(z, s)$ if $\sigma \cdot \infty = \infty$.

q.e.d.

COROLLARY. *For any given i^* in k we have*

$$E(z, s) - |v|_A^s |N(z + i^*)|_A^{-s} < |v|_A^{s-1}$$

on the subset of \mathcal{X} defined by

$$(*) \quad |N(z + i^*)|_A < |v|_A^2 < 1;$$

we also have

$$E(z, s) - |v|_A^s |N(z+i^*)|_A^{-s} < 1$$

on the subset of \mathcal{X} defined by

$$(**) \quad |N(z+i^*)|_A < |v|_A .$$

The proof is quite simple: choose an element σ of $GL_2(k)$ satisfying $\sigma^{-1} \cdot \infty = -i^*$; then by Th. 2 we get

$$E(z, s) - |v|_A^s |N(z+i^*)|_A^{-s} = E(\sigma \cdot z, s) - |v(\sigma \cdot z)|_A^s .$$

And (*) implies

$$|v(\sigma \cdot z)|_A = |v|_A |N(z+i^*)|_A^{-1} > |v|_A^{-1} > 1$$

while (**) implies $|v(\sigma \cdot z)|_A > 1$. The rest follows from Th. 1.

One way to construct a subset of \mathcal{X} satisfying (*) is as follows: we restrict v as $|v|_A < 1$ and u by the condition that $(u+i^*)v^{-1}$ remains in a compact subset of k_A , i.e., as $|u+i^*|_p < |v|_p$ for every p in a fixed finite set and $|u+i^*|_p \leq |v|_p$ for every other p . In fact by Lemma 2 we then get

$$\begin{aligned} |N(z+i^*)|_A &> \prod_p \max(|u+i^*|_p, |v|_p)^2 \\ &> \prod_p |v|_p^2 = |v|_A^2 < 1 . \end{aligned}$$

We shall examine our definitions and results in the special case where $k = \mathbf{Q}$. In order to make the objects visible, we shall assume that $u_0 = 0$ and $v_0 = 1$ in $z = u + vw = (z_0, z_\infty)$; then we get

$$E(z, s) = |\text{Im}(z_\infty)|^s (1 + \sum_{\gamma, \delta} |\gamma z_\infty + \delta|^{-2s}) ,$$

in which γ, δ are relatively prime integers with $\gamma \geq 1$. This is a classical Eisenstein series; cf., e.g., Kubota [5]. Moreover the condition (*) and the more restrictive condition explained after the corollary both become

$$|z_\infty + i^*| < |\text{Im}(z_\infty)| < 1 ;$$

if we restrict z_∞ to the upper-half plane, this defines a ∇ -shaped region which touches the real axis at $-i^*$. On the other hand the condition (**) becomes

$$|z_\infty + i^*|^2 < |\text{Im}(z_\infty)| ;$$

this defines a closed circle which is tangent to the real axis at $-i^*$. In the first case the height and the angle between the two lines through $-i^*$ are arbitrary and in the second case the radius is arbitrary. We have thus shown that

the second part of Th. 1 and the corollary become the well-known theorem describing the behavior of the classical Eisenstein series at cusps; cf. Kubota, op. cit.

4. **Main results.** We shall go back to our $(E' - E)(U(g)\Phi)$: in view of Th. 1 we put

$$|g|_A = |(u, t)|_A = |t|_A^m, \quad s = n/2m;$$

then $s > 1$ and $g \rightarrow |g|_A$ gives a continuous homomorphism of G_A to \mathbf{R}_+^* .

THEOREM 3. *We restrict g as $|g|_A > 1$ and Φ to a compact subset of $\mathcal{S}(X_A)$; then we have*

$$(E' - E)(U(g)\Phi) < |g|_A^{1-s}.$$

The proof of this theorem does not depend on the degree m of $f(x)$; it can be extracted from the proof of Prop. 7 in Weil [12], Chap. V. Since we have included the function field case, we shall give an outline:

LEMMA 3. *Let F denote a closed subset of $X_A = X_0 \times X_\infty$ such that $F \cap (X_0 \times \{0\}) = \emptyset$; then for any $N \geq 0$ and any compact subset C of $\mathcal{S}(S_A)$ there exists an element ϕ of $\mathcal{S}(X_A)$ satisfying*

$$|a_\tau|_A^N |\Phi(a_\tau x)| \leq \phi(x)$$

for every x in F , every a_τ with $|a_\tau|_A \geq 1$, and every Φ in C .

PROOF. This lemma is known if k is a number field; cf. Weil [12], Chap. I, Lemma 6 and also Mars [8], Lemma 7. The proof in the function field case is simpler: by using coordinates of x_∞ we define $|x_\infty|_\infty$ as

$$|x_\infty|_\infty = \max_{1 \leq i \leq n} \{|x_{\infty, i}|_{p_\infty}\}.$$

Since C is compact, there exist compact open subgroups L_0, L_∞ of X_0, X_∞ , respectively, such that $\text{Supp}(\Phi)$ is contained in $L_0 \times L_\infty$ for every Φ in C ; we may assume that L_∞ is defined by $|x_\infty|_\infty \leq c$ for some $c > 0$. Let ϕ_0 denote a constant multiple of the characteristic function of L_0 and ϕ_∞ the characteristic function of L_∞ ; then if we take the constant large enough, we will have

$$|\Phi(x)| \leq \phi_0(x_0)\phi_\infty(x_\infty)$$

for every $x = (x_0, x_\infty)$ in X_A and Φ in C . Since $F \cap (L_0 \times L_\infty)$ is compact and has an empty intersection with $X_0 \times \{0\}$, we have $|x_\infty|_\infty \geq d$ for some $d > 0$ and for every x in $F \cap (L_0 \times L_\infty)$. We shall show that

$$\phi = (cd^{-1})^N (\phi_0 \otimes \phi_\infty)$$

has the required property.

If $\Phi(a_\tau x) \neq 0$ for some x in X_A and Φ in C , we get $\phi_0(x_0)\phi_\infty(\tau x_\infty) \neq 0$; hence x_0 is in L_0 and τx_∞ is in L_∞ . If $|a_\tau|_A = |\tau|_\infty \geq 1$, then x_∞ itself is in L_∞ . And if x is also in F , then $|\tau|_\infty \leq c|x_\infty|_\infty^{-1} \leq cd^{-1}$; hence

$$\begin{aligned} |a_\tau|_A^N |\Phi(a_\tau x)| &\leq (cd^{-1})^N \phi_0(x_0)\phi_\infty(\tau x_\infty) \\ &= (cd^{-1})^N \phi_0(x_0)\phi_\infty(x_\infty). \end{aligned} \quad \text{q.e.d.}$$

The proof of Th. 3 is as follows: we have shown in §1 that

$$G_A = G_k \cdot \{(u, a_\tau); u_0 = 0, u_\infty \in \text{compact}\} \cdot \text{compact}.$$

If we express an element g of G_A as the product of an element of G_k , (u, a_τ) , and an element of the compact set, then we have $|g|_A < |a_\tau|_A^m$. Therefore we may assume that $g = (u, a_\tau)$, in which $u_0 = 0$, u_∞ is in a fixed compact subset of k_∞ , and $|a_\tau|_A > 1$; then $(a_\tau^{-m}u, 1)$ in

$$(u, a_\tau) = (0, a_\tau)(a_\tau^{-m}u, 1)$$

remains in a compact subset of G_A . Therefore we may assume that $g = (0, a_\tau)$ where $|a_\tau|_A \geq 1$. We recall that for $g = (0, t)$ we have

$$(E' - E)(U(g)\Phi) = |t|_A^{(1/2)n} \left(\sum_{\xi \in X_k - \{0\}} \Phi(t\xi) - \sum_{i \in k} \int_{U(i)_A} \Phi(tx)|\theta_i(x)|_A \right).$$

We shall take $t = a_\tau$ where $|a_\tau|_A \geq 1$. Then by applying Lemma 3 to $F = X_k - \{0\}$ we get

$$|t|_A^{(1/2)n} \cdot \sum_{\xi \in X_k - \{0\}} \Phi(t\xi) < |g|_A^{-N}$$

for any $N \geq 0$; by applying Lemma 3 to $F = f^{-1}(k^\times)$ we get

$$|t|_A^{(1/2)n} \cdot \sum_{i \in k^\times} \int_{U(i)_A} \Phi(tx)|\theta_i(x)|_A < |g|_A^{-N}$$

for any $N \geq 0$; and finally we have

$$|t|_A^{(1/2)n} \cdot \int_{U(0)_A} \Phi(tx)|\theta_0(x)|_A = |g|_A^{-s} \cdot \int_{U(0)_A} \Phi|\theta_0|_A.$$

We have only to put the above three relations together.

THEOREM 4. *Let i^* denote any given element of k ; we restrict $g = (u, t)$ by $|t|_A < 1$ and by the additional condition that $(u + i^*)t^{-m}$ remains in a compact*

subset of k_A ; we also restrict Φ to a compact subset of $\mathcal{S}(X_A)$. Then we have

$$(E' - E)(U(g)\Phi) < |g|_A^{s-1}.$$

PROOF. We define the singular term at $-i^*$ as

$$\text{s.t. at } -i^* = |t|_A^{(1/2)n} \cdot \int_{X_A} \Phi((u+i^*)f(x))\Phi(tx)|dx|_A;$$

then by Th. 1, its corollary, and a remark after the corollary we get

$$E(U(g)\Phi) - \text{s.t. at } -i^* < |g|_A^{s-1}$$

on the subset under consideration. Therefore we have only to show that

$$E'(U(g)\Phi) - \text{s.t. at } -i^* < |g|_A^N$$

for any $N \geq 0$. Since $\phi(uf(\xi)) = \phi((u+i^*)f(\xi))$ for every ξ in X_k , by applying the usual Poisson formula we get

$$\begin{aligned} E'(U(g)\Phi) &= |t|_A^{(1/2)n} \cdot \sum_{\xi \in X_k} \phi((u+i^*)f(\xi))\Phi(t\xi) \\ &= |t|_A^{(1/2)n} \cdot \sum_{\xi \in X_k} \int_{X_A} \phi([\xi, y] + (u+i^*)f(y))\Phi(ty)|dy|_A. \end{aligned}$$

Therefore if we put

$$\Psi(x) = \phi((u+i^*)t^{-m}f(x))\Phi(x)$$

and denote its Fourier transform by $\Psi^*(x)$, we get

$$E'(U(g)\Phi) - \text{s.t. at } -i^* = |t|_A^{-(1/2)n} \cdot \sum_{\xi \in X_k - \{0\}} \Psi^*(t^{-1}\xi).$$

We recall that the mapping $G_A \times \mathcal{S}(X_A) \rightarrow \mathcal{S}(X_A)$ defined by $(g, \Phi) \rightarrow U(g)\Phi$ is continuous. Since $(u+i^*)t^{-m}$ remains in a compact subset of k_A , therefore, the set $\{\Psi\}$ is relatively compact in $\mathcal{S}(X_A)$. Since the Fourier transformation is bicontinuous, the set $\{\Psi^*\}$ is also relatively compact in $\mathcal{S}(X_A)$. The remaining part of the proof is as follows:

We write $t = ia_\tau c$ with i in k^\times and c in a fixed compact subset of k_A^\times ; then by assumption we have $|t|_A \asymp |a_\tau|_A < 1$ and

$$|t|_A^{-(1/2)n} \cdot \sum_{\xi \in X_k - \{0\}} \Psi^*(t^{-1}\xi) \asymp |a_\tau^{-1}|_A^{(1/2)n} \cdot \sum_{\xi \in X_k - \{0\}} \Psi^*(c^{-1}a_\tau^{-1}\xi).$$

Since the set $\{\Psi^*(c^{-1}x)\}$ with c in the compact subset is still relatively compact and $|a_\tau^{-1}|_A > 1$, we can apply Lemma 3 to $F = X_k - \{0\}$. In this way we get

$$|a_\tau^{-1}|_A^{(1/2)n} \cdot \sum_{\xi \in X_k - \{0\}} \Psi^*(c^{-1}a_\tau^{-1}\xi) < |a_\tau^{-1}|_A^{mN} \asymp |g|_A^N$$

for any $N \geq 0$.

q.e.d.

We shall examine Th. 4 in the special case where $k=Q$; for the sake of simplicity we shall assume that $f(x)$ has integer coefficients: we shall use the particular ϕ defined by the condition that $\phi_\infty(u_\infty)=e(u_\infty)$ for every u_∞ in $\mathbf{Q}_\infty=\mathbf{R}$ and $\phi_p(u_p)=\phi_\infty(-\langle u_p \rangle)$ for every u_p in the Hensel p -adic field \mathbf{Q}_p , where $\langle u_p \rangle$ is the fractional part of u_p . We shall use as Φ a function of the form $\Phi_0 \otimes \Phi_\infty$, in which Φ_∞ is a Schwartz function on $X_\infty=\mathbf{R}^n$; for the sake of simplicity we shall assume that Φ_0 is the characteristic function of X_0° . And we shall take as g a special pair (u, a_τ) where $u_0=0$. Then we get

$$(U(g)\Phi)(x)=\tau^{(1/2)n}e(u_\infty f(x_\infty))\Phi_\infty(\tau x_\infty)$$

if x_0 is in X_0° and $(U(g)\Phi)(x)=0$ otherwise; hence

$$(1) \quad \tau^{-(1/2)n}E'(U(g)\Phi)=\sum_{\xi \in \mathbf{Z}^n} e(u_\infty f(\xi))\Phi_\infty(\tau\xi).$$

One way to visualize this series is to consider the limit case where Φ_∞ becomes the characteristic function of a relatively compact, say open, subset J of \mathbf{R}^n ; then we get the following finite sum:

$$\sum_{\xi \in \mathbf{Z}^n \cap \tau^{-1}J} e(u_\infty f(\xi)).$$

This is classically known as an exponential sum; it behaves quite delicately as a function of u_∞ and τ ; cf., e.g., Birch [2].

We shall also make the two series for $E(U(g)\Phi)$ explicit: suppose that $i^*=\gamma^{-1}\delta$, where γ, δ are relatively prime integers with $\gamma \geq 1$, and put

$$G(i^*)=\gamma^{-n} \cdot \sum_{\xi \bmod \gamma} e(-i^* f(\xi));$$

then one expression for $\tau^{-(1/2)n}E(U(g)\Phi)$ is

$$(2) \quad 1 + \sum_{i^* \in Q} G(i^*) \cdot \int_{\mathbf{R}^n} e((u_\infty + i^*)f(x_\infty))\Phi_\infty(\tau x_\infty) dx_\infty,$$

in which dx_∞ is the usual measure on \mathbf{R}^n . On the other hand, for every integer i and a positive integer Q let $N_Q(i)$ denote the number of $\xi \bmod Q$ satisfying $f(\xi) \equiv i \bmod Q$; we let Q tend to ∞ so that Q becomes divisible by any positive integer; then the following limit:

$$S(i) = \lim_{Q \rightarrow \infty} Q^{-(n-1)} N_Q(i)$$

exists. In fact the condition of convergence is $n \geq 4$ if $i \neq 0$ and $n > \max(m+1, 4)$ if $i=0$; and this is weaker than our assumption that $n > 2m$ (and $m \geq 2$). We might also mention that in the fourth and fifth lines of [4], p. 224 the condition

$n > \max(m+1, 4)$ was incorrectly stated as $n > \max(m, 4)$. At any rate another expression for $\tau^{-(1/2)n} E(U(g)\Phi)$ is

$$(3) \quad 1 + \sum_{i \in \mathbb{Z}} S(i) \cdot \int_{f^{-1}(i)} \Phi_{\infty}(\tau x_{\infty}) | \theta_i(x_{\infty}) |_{\infty} \cdot e(iu_{\infty}) .$$

And Th. 4 shows that as the point

$$z_{\infty} = u_{\infty} + (-1)^{1/2} \tau^m$$

approaches any rational number $-i^*$ in the ∇ -shaped region explained in the previous section, (1) is given by (2) and (3) with a remainder term of order τ^{-m} ; we observe that (1) itself can be of order τ^{-n} .

We shall also explain the classical case where $f(x)$ is a quadratic form: after Hermite and Siegel we choose a "majorant" $h(x_{\infty})$ of $f(x_{\infty})$ and put

$$\Phi_{\infty}(x_{\infty}) = \exp(-2\pi h(x_{\infty}));$$

we can define $h(x_{\infty})$ as a positive-definite quadratic form (with real coefficients) such that the above Φ_{∞} is equal to its Fourier transform relative to the bicharacter $\phi_{\infty}(f(x_{\infty}, y_{\infty}))$ of $X_{\infty} \times X_{\infty}$, in which

$$f(x, y) = f(x+y) - f(x) - f(y) .$$

If $f(x_{\infty})$ is positive-definite, then $h(x_{\infty}) = f(x_{\infty})$ is the only choice. At any rate under the above specialization (1) becomes the following theta series:

$$\sum_{\xi \in \mathbb{Z}^n} e(\operatorname{Re}(z_{\infty})f(\xi) + (-1)^{1/2} \operatorname{Im}(z_{\infty})h(\xi))$$

and (2) becomes the following Eisenstein series:

$$1 + e((1/8)(p-q)) |d|^{-(1/2)} \cdot \sum_{i^* \in \mathcal{Q}} G(i^*) (z_{\infty} + i^*)^{-(1/2)p} (\bar{z}_{\infty} + i^*)^{-(1/2)q} ,$$

in which p, q are the numbers of positive and negative eigenvalues of the coefficient matrix of $f(x, y)$ and d is its determinant. Furthermore in the special case where $q=0$ (3) becomes

$$1 + (2\pi)^{(1/2)n} \Gamma\left(\frac{1}{2}n\right)^{-1} d^{-(1/2)} \cdot \sum_{i=1}^{\infty} S(i) i^{(1/2)n-1} e(iz_{\infty}) .$$

We recall that such series and their generalizations appeared in many works of Siegel.

5. A conjecture. In the special case where $f(x)$ is a quadratic form we know that $(E' - E)(U(g)\Phi)$ is bounded on G_A if Φ is restricted to a compact subset of $\mathcal{S}(X_A)$. We refer to Weil [12], Prop. 7 for the proof; we also refer to Siegel

[9] and Arıturk [1]. We have shown in the general case that $(E' - E)(U(g)\Phi)$ vanishes to the order $s-1$ at every k -rational boundary point of G_A and hence is bounded around every such point. Therefore it is reasonable to think that the boundedness is not restricted to the case of a quadratic form. We in fact propose the following weaker statement as a conjecture:

CONJECTURE. For any given Φ in $\mathcal{S}(X_A)$ there exists a positive real number ε such that

$$(E' - E)(U(g)\Phi) < |g|_A^{1-s+\varepsilon}$$

or at least

$$\int_{k_A/k} (E' - E)(U(g)\Phi) |du|_A < |g|_A^{1-s+\varepsilon}$$

on the subset of G_A defined by $|g|_A \leq 1$.

We observe that $(E' - E)(U(u, t)\Phi)$ is a continuous function on $(k_A/k) \times k_A^\times$ and that the first hypothesis implies the second. We can prove this conjecture if n is sufficiently large compared to m . The point is that it may be true under the assumption $n > 2m$ and it is very likely to be true if $n > m^2$ and that it implies the following generalization of the Hasse-Minkowski theorem:

"If the non-singular projective hypersurface defined by $f(x) = 0$ has a k_p -rational point for every p , then it has a k -rational point".

(In this way we can reproduce Birch's result in [2] at least for a non-singular projective hypersurface.) The above implication can be proved as follows: we have

$$(E' - E)(\mathbf{t}(u)\Phi) = \sum_{i \in k} c_i(\Phi) \phi(iu),$$

in which

$$c_i(\Phi) = \sum_{\xi \in U^{(i)}_k} \Phi(\xi) - \int_{U^{(i)}_A} \Phi | \theta_i |_A$$

for every i in k . Since we also have

$$c_i(\Phi) = \int_{k_A/k} (E' - E)(\mathbf{t}(u)\Phi) \phi(-iu) |du|_A,$$

if the conjecture is true, we will have

$$c_0(\mathbf{d}(t)\Phi) < |t|_A^{m(1-s+\varepsilon)}$$

on the subset of k_A^\times defined by $|t|_A \leq 1$. By using the O -symbol as $|t|_A \rightarrow 0$, this can be rewritten as

$$\sum_{\xi \in U(0)_k} \Phi(t\xi) = \left(\int_{U(0)_A} \Phi|\theta_0|_A + O(|t|_A^{m\epsilon}) \right) |t|_A^{m-n}.$$

We recall that for every i in k the support of the measure $|\theta_i|_A$ is the whole space $U(i)_A$. Therefore if $U(0)_A$ is not empty, for any $\Phi \geq 0$ which is not the constant 0 on $U(0)_A$ its integral over $U(0)_A$ is positive; then the right hand side, hence also the left hand side, is positive for all small $|t|_A$. Therefore $U(0)_k$ can not possibly be empty; this completes the proof.

Appendix

One of the points we have made in this paper is that we can do something significant by using only such a small group as G_A . However the last section clearly indicates that a good generalization of the metaplectic group will be extremely useful. For that purpose we must start examining, as in Kubota [6], possible generalizations of the metaplectic group over a local field. Since the problem is still quite difficult, it was suggested to us by Shalika to examine the same problem over a finite field. In this appendix we shall explain our fragmental results in this simplest case.

We shall introduce a naive generalization of the metaplectic group over $K = F_q$; let X denote an n -dimensional vector space over K and $L(X)$ the Hilbert space of (complex-valued) functions on X with the following norm:

$$\|\Phi\|^2 = \text{card}(X)^{-(1/2)} \cdot \sum_{x \in X} |\Phi(x)|^2;$$

then for every t in K^\times the scalar multiplication by t in X defines a unitary operator in $L(X)$. Let ϕ denote a non-trivial character of K and $f(x)$ a form of degree m on X , i.e., a homogeneous element of degree m of the symmetric algebra of the dual of X ; then for every u in K the multiplication by $\phi(uf(x))$ in $L(X)$ also defines a unitary operator in $L(X)$. Finally let $[x, y]$ denote a symmetric non-degenerate K -bilinear form on $X \times X$; then

$$\Phi^*(x) = \text{card}(X)^{-(1/2)} \cdot \sum_{y \in X} \phi([x, y]) \Phi(y)$$

defines a unitary operator in $L(X)$. We shall assume that $f(x)$ is non-degenerate, i.e., $f(x)$ is not a form on the quotient of X by a proper subspace; and we shall denote by Mp the subgroup of the full unitary group $\text{Aut}(L(X))$ generated by the above three types of unitary operators.

The structure of Mp is known for $m \leq 2$: if $m=1$, we necessarily get $n=1$; and Mp becomes a semidirect product of a Heisenberg group of order q^3 by a

generalized quaternion group of order $2(q-1)$. If $m=2$, q is odd, and $[x, y]=f(x, y)$, then Mp becomes a semidirect product of a cyclic group of order 2 or 4 (according as $q \equiv \pm 1 \pmod 4$) by $SL_2(K)$. In the general case we have the following theorem:

THEOREM 5. *Suppose that $n=1$ and let e denote the G.C.D. of m and $q-1$. Then if e is even, Mp decomposes into $\frac{1}{2}e+1$ inequivalent irreducible representations of degrees*

$$\frac{q-1}{e}+1, \frac{q-1}{e}, \frac{2(q-1)}{e}, \dots, \frac{2(q-1)}{e};$$

and if e is odd, Mp decomposes into $\frac{1}{2}(e+1)$ inequivalent irreducible representations of degrees

$$\frac{q-1}{e}+1, \frac{2(q-1)}{e}, \dots, \frac{2(q-1)}{e}.$$

PROOF. For a moment we shall drop the assumption that $n=1$; let T denote an element of $\text{End}(L(X))$, i.e., a K -linear transformation of $L(X)$ to itself; let δ_y denote the Dirac function on X satisfying $\delta_y(y)=1$ and put $k(x, y)=(T\delta_y)(x)$; then we get

$$(T\Phi)(x) = \sum_{y \in X} k(x, y)\Phi(y)$$

for every Φ in $L(X)$. The set A of all T 's which elementwise commute with Mp forms a subalgebra of $\text{End}(L(X))$. Let $k(x, y)$ denote the kernel of T defined as above; then T is in A if and only if

- (k1) $k(tx, ty)=k(x, y)$ for every t in K^\times ;
- (k2) $k(x, y) \neq 0$ only if $f(x)=f(y)$;
- (k3) $\sum_{z \in X} k(x, z)\phi([z, y]) = \sum_{z \in X} \phi([x, z])k(z, y)$.

We shall closely examine these conditions resuming the assumption that $n=1$: if we identify X with K , we get $f(x)=cx^m$, $[x, y]=dxy$ for some c, d in K^\times . We define a function ϕ on K^\times as $\phi(t)=k(1, t)$; then we get

- (\phi1) $\text{Supp}(\phi) \subset (K^\times)_e$, i.e., $\phi(t) \neq 0$ only if $t^e=1$;
- (\phi2) $\phi(t)=\phi(t^{-1})$.

Conversely for every such ϕ if we define $k(x, y)$ as

$$k(x, y) = \begin{cases} \phi(x^{-1}y) & xy \neq 0 \\ 0 & xy = 0, \quad (x, y) \neq (0, 0) \\ \sum_t \phi(t) & (x, y) = (0, 0), \end{cases}$$

then k satisfies (k1)-(k3). And the correspondence $k \rightarrow \phi$ gives rise to a C -linear bijection.

Consider the group ring $C(K^\times)$ of K^\times ; then the correspondence

$$T \rightarrow k \rightarrow a = \sum_t \phi(t)t$$

gives rise to an injective algebra homomorphism of A to $C(K^\times)$. In particular A is commutative. Since Mp is completely reducible, therefore, every irreducible subrepresentation of Mp has multiplicity one. As for its degree, it can be determined as follows:

We observe that the C -linear extension, say ω , of a character χ of K^\times is an algebra homomorphism $C(K^\times) \rightarrow C$ mapping 1 to 1 and the correspondence $\chi \rightarrow \omega$ gives rise to a bijection. In particular there are $q-1$ such algebra homomorphisms. If we arrange them in some order, we get a representation of $C(K^\times)$ as the algebra of all diagonal matrices of degree $q-1$; and by restricting this representation to A we get a diagonalization of A .

After this remark suppose that ω, ω' correspond to χ, χ' , respectively; then we get $\omega' = \omega$ on A if and only if $\chi' = \chi^{\pm 1}$ on $(K^\times)_e$. Therefore the restriction of ω to A has multiplicity $(q-1)/e$ or $2(q-1)/e$ according as $\chi^2 = 1$ or $\chi^2 \neq 1$ on $(K^\times)_e$. Moreover if ω_0 corresponds to the character 1 of K^\times , then we get

$$k(0, 0) = \sum_t \phi(t) = \omega_0(a).$$

We have thus obtained a diagonalization of A as an algebra of q -by- q matrices from which we can read off the degrees of irreducible subrepresentations of Mp . q.e.d.

COROLLARY. *If q is odd, $e \geq 3$, and $(e, q) \neq (3, 7), (4, 5)$, then Mp is not a central extension of $SL_2(K)$; if q is odd and $e \geq 5$, then every homomorphism of $SL_2(K)$ to Mp is trivial.*

The proof is as follows: we recall that the degrees of irreducible representations of $SL_2(K)$ are

$$1, \quad \frac{1}{2}(q \pm 1), \quad q \pm 1, \quad q;$$

cf. Tanaka [10]. And we have

$$\frac{1}{2}(q-1) > 2(q-1)/e$$

if (and only if) $e \geq 5$. Therefore if $e \geq 5$, the degree of any non-trivial irreducible

representation of $SL_2(K)$ is larger than the degrees of irreducible subrepresentations of Mp . After this observation suppose that π is a homomorphism of $SL_2(K)$ to Mp ; then π followed by any irreducible subrepresentation of Mp is necessarily a sum of the trivial representation of $SL_2(K)$. Therefore π is trivial; this proves the second part and also the first part for $e \geq 5$ (because any central extension of $SL_2(K)$ splits). The proof of the remaining cases is similar.

References

- [1] Ariturk, H.: The Siegel-Weil formula for orthogonal groups, Thesis, Johns Hopkins University, 1975.
- [2] Birch, B. J.: Forms in many variables, Proc. Royal Soc. A, **265** (1962), 245-263.
- [3] Igusa, J.: On the arithmetic of Pfaffians, Nagoya Math. J. **47** (1972), 169-198.
- [4] Igusa, J.: On a certain Poisson formula, Nagoya Math. J. **53** (1974), 211-233.
- [5] Kubota, T.: Elementary theory of Eisenstein series, Kodansha, Tokyo, 1973.
- [6] Kubota, T.: A generalized Weil type representation, TR 73-7, University of Maryland, 1973.
- [7] Levin, M.: A continuity problem in the Siegel-Weil formula, TR 74-10, University of Maryland, 1974.
- [8] Mars, J. G. M.: Les nombres de Tamagawa de certains groupes exceptionnels, Bull. Soc. Math. France, **94** (1966), 97-140.
- [9] Siegel, C. L.: Gesammelte Abhandlungen, I-III, Springer, 1966; in particular, Indefinite quadratische Formen und Funktionentheorie I, Math. Ann. **124** (1951), 17-54; III, 105-142.
- [10] Tanaka, S.: Construction and classification of irreducible representations of special linear group of the second order over a finite field, Osaka J. Math. **4** (1967), 65-84.
- [11] Tate, J.: Fourier analysis in number fields and Hecke's zeta-functions, Thesis, Princeton University, 1950; Algebraic number theory, Thompson Book Co., 1967, 305-347.
- [12] Weil, A.: Sur la formule de Siegel dans la théorie des groupes classiques, Acta Math. **113** (1965), 1-87.
- [13] Weil, A.: Basic number theory, Grundle. Math. Wiss. 144, Springer, 1967.

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