# SPECTRAL DECOMPOSITION OF $L^{2}(N \backslash \mathrm{GL}(2), \eta)$ 

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#### Abstract

Let $G$ denote $\operatorname{GL}(2, F)$, where $F$ is $\mathbf{R}$ or $\mathbf{C}$ or a $p$-adic field. Let $\eta$ be a non-trivial character of the unipotent upper triangular group $N$ in $G$. The object of this paper is to present an explicit spectral decomposition of $L^{2}(N \backslash G, \eta)$, the representation of $G$ induced unitarily by $\eta$. This representation is well-known to be multiplicity-free and to be quasiequivalent to the (right) regular representation of $G$.


Introduction. In §1 we give a 'cuspidal' characterization of the discrete spectrum of $L^{2}(N \backslash \operatorname{SL}(2, F), \eta)$. In $\S 2$ we prove a crucial duality formula, which later allows us to decompose the continuous spectrum. It also 'explains' the occurrence of $\varepsilon^{\prime}$-factors in the measure giving the direct integral decomposition.

The scalar product of $\S 3$ is originally due, in the $p$-adic case, to $H$. Jacquet. Here we suitably modify and extend his (unpublished) work and also treat the archimedean case. We have tried to present a unified approach.

This work was done as a part of my doctoral thesis [11] at Columbia. Besides the obvious debt to my advisor then, H. Jacquet, I would like to mention the strong influence of R. Godement's paper ([3]) on the spectral analysis of modular functions, and thank P. Sally for his interest and critical comments. I would like to thank the referee for his helpful remarks which led to simplifications of some of the proofs. Thanks are also due to Miss M. Murray and Mrs. Anne Wolfsheimer for their excellent typing of this manuscript.
0. Terminology. Fix a non-trivial character $\psi$ of $F^{+}$. Let $d x$ be the self-dual measure on $F$ with respect to $\psi$ and let | | be the normalized absolute value. When $F$ is non-archimedean, let $\mathfrak{D}, \mathfrak{p}, v, \pi$ and $\mathbf{F}_{q}$ respectively denote the ring of integers, maximal ideal, valuation, uniformizer and residue field. On $F^{2}$ we will always take the product measure induced by the self-dual measure $d x$ on $F$. We will denote by $U_{F}$ the group $\{x \in F||x|=1\}$.

Let

$$
\begin{aligned}
& G=\mathrm{GL}(2, F), \quad N=\left\{\left(\begin{array}{cc}
1 & * \\
0 & 1
\end{array}\right)\right\} \subset G \\
& A=\left\{\left.\left(\begin{array}{cc}
a_{1} & 0 \\
0 & a_{2}
\end{array}\right) \right\rvert\, a_{1}, a_{2} \in F^{*}\right\}, \quad B=N A, \\
& Z=\left\{\left.\left(\begin{array}{ll}
t & 0 \\
0 & t
\end{array}\right) \right\rvert\, t \in F^{*}\right\}, \quad P=\left\{\left(\begin{array}{ll}
* & * \\
0 & 1
\end{array}\right)\right\}, \quad \text { and } \\
& K= \begin{cases}\mathrm{GL}_{2} \mathcal{D}, & \text { if } F \text { is } p \text {-adic } \\
O(2), & \text { if } F=\mathbf{R} \\
U(2), & \text { if } F=\mathbf{C} .\end{cases}
\end{aligned}
$$

Denote by $\eta$ the character of $N$ given by: $\left(\begin{array}{ll}1 & x \\ 0 & 1\end{array}\right) \mapsto \psi(x)$. Then $L^{2}(N \backslash G, \eta)$ consists of (classes of) measurable functions $f$ on $G$ which are square-integrable $\bmod N$ and satisfy:

$$
f(n g)=\eta(n) f(g), \quad \text { for almost all }(n, g) \text { in } N \times G
$$

Let $G^{1}=\operatorname{SL}(2, F)$. For every subgroup $H$ of $G$, let $H^{1}$ (resp. $H_{K}$ ) denote $H \cap G^{1}$ (resp. $H \cap K$ ).

Let $\delta: A \rightarrow \mathbf{R}_{+}^{*}$ be the module $A$ defined by: $\delta(a)=\left|a_{1} / a_{2}\right|$, for

$$
a=\left(\begin{array}{cc}
a_{1} & 0 \\
0 & a_{2}
\end{array}\right)
$$

Let $d^{*} a$ denote the product measure on $A \stackrel{\approx}{\rightarrow}\left(F^{*}\right)^{2}, a \mapsto\left(a_{1}, a_{2}\right)$, so that $d^{*} a=d^{*} a_{1} \cdot d^{*} a_{2}$.

Fix a Haar measure $d k$ on $K$ such that $d k$ coincides with $d^{*} a$ on $A_{K}$. Then the quotient measure on $N \backslash G$ may be identified, via Iwasawa decomposition: $G=N A K$, with: $\delta(a)^{-1} \cdot d^{*} a d k$.

If $h_{1}, h_{2}$ are in $L^{2}(N \backslash G, \eta)$ (resp. $L^{2}(N \backslash G)$ ), the function: $g \mapsto$ $h_{1}(g) \bar{h}_{2}(g)$ is invariant on the left by $N$. We set

$$
\left\langle h_{1}, h_{2}\right\rangle=\int_{N \backslash G} h_{1}(g) \overline{h_{2}(g)} d g
$$

For any locally compact topological group $H$, let $\hat{H}$ denote the Mackey dual of $H$, viz. the set of classes of inequivalent, irreducible, unitary representations of $H$. By abuse of notation, we will often use the same symbol for an element of $\hat{H}$ and for a representative of that element.

If $S$ is a topological space, and if $f_{1}, f_{2}$ are two complex-valued functions on $S$, then we write: $f_{1} \prec f_{2}$ if there exists a positive scalar $C$
such that $\left|f_{1}(x)\right| \leq C\left|f_{2}(x)\right|$, for all $x$ in $S$. Sometimes we will also write: $f_{1}(x) \prec f_{2}(x)$ for all $x \in S$. In particular, $f_{1} \prec 1$ implies that $f_{1}$ is bounded on $S$.

We write $f_{1} \succ f_{2}$ if $f_{2} \prec f_{1}$, and $f_{1} \asymp f_{2}$ if we have simultaneously: $f_{1} \prec f_{2}$ and $f_{1} \succ f_{2}$.

We will always write, for a positive real number $t$, that $t \ll 1$ (resp. $t \gg 1$ ) if $t$ is very close to $0($ resp. to $+\infty)$.

1. The intertwining operator $\theta$. We first recall the natural identification of the homogeneous space $N \backslash G^{1}$ with $F^{2}-\{0\}$.
$G^{1}$ acts by right multiplication $\rho$ on $F^{2}=\{(x, y) \mid x, y \in F\}$. This action leaves $F^{2}-\{0\}$ stable and is transitive on it. Furthermore, if $e_{1}$ denotes the vector $(0,1)$ in $F^{2}$, then $N$ is precisely the stabilizer of $e_{1}$ in $G^{1}$. This results in a right $G^{1}$-space analytic isomorphism:

$$
\begin{gather*}
N \backslash G^{1} \xrightarrow[\rightarrow]{\sim} F^{2}-\{0\}  \tag{1.0}\\
N g \mapsto \rho_{g} \cdot e_{1}=e_{1} \cdot g
\end{gather*}
$$

Let $S\left(F^{2}\right)$ denote the space (cf. [7]) of Schwartz-Bruhat functions on $F^{2}$. It is stable under the right action $\rho$ of $G^{1}$.

Definition 1.1. For $\Phi$ in $S\left(F^{2}\right)$ and $g$ in $G^{1}$, set:

$$
f_{\Phi}(g)=\Phi\left(e_{1} \cdot g\right)
$$

Since the Haar measure on $F^{2}$ is invariant under the action of $G^{1}$, it must be, under the above identification, the invariant measure on $N \backslash G^{1}$. Hence we obtain:

Lemma 1.2. $f_{\Phi} \in L^{2}\left(N \backslash G^{1}\right)$, for every $\Phi \in S\left(F^{2}\right)$.
Definition 1.3.

$$
S\left(N \backslash G^{1}\right)=\left\{f_{\Phi} \mid \Phi \in S\left(F^{2}\right)\right\}
$$

Note that the map $\Phi \mapsto f_{\Phi}$ is injective and is equivariant with respect to the right action of $G^{1}$.

The following proposition is an immediate consequence of the remark preceding Lemma 1.2.

Proposition 1.4. The $G^{1}$-equivariant isomorphism

$$
S\left(F^{2}\right) \stackrel{\sim}{\rightarrow} S\left(N \backslash G^{1}\right), \quad \Phi \mapsto f_{\Phi}
$$

extends to a Hilbert space isomorphism of the unitary $G^{1}$-modules $L^{2}\left(F^{2}\right)$ and $L^{2}\left(N \backslash G^{1}\right)$.

Definition 1.5. Let $\mathscr{L}$ denote the set of all lines in $F^{2}$ passing through the origin, and set:

$$
\begin{aligned}
S_{0}\left(F^{2}\right)=\{ & \left\{\in S\left(F^{2}\right) \mid \int_{L} \Phi(z) d z=0\right. \\
& \forall \text { line } L \text { in } \mathscr{L} \text { with Haar measure } d z\} .
\end{aligned}
$$

Proposition 1.6. $S_{0}\left(F^{2}\right)$ is $G^{1}$-stable, and is dense in $L^{2}\left(F^{2}\right)$.
Proof. The right action of $G^{1}$ on $F^{2}$ leads to an action of $G^{1}$ on the space $\mathscr{L}$ of all lines in $F^{2}$ as follows: Given $L$ in $\mathscr{L}$, choose $w$ in $F^{2}$ such that $L=F w$. If $g$ is in $G^{1}$, set: $\rho_{g}(L)=F \cdot\left(\rho_{g} w\right)$. It may be checked that $\rho_{g}(L)$ is independent of the choice of $w$.

For every $\Phi$ in $S\left(F^{2}\right)$ and for every $L$ in $\mathscr{L}$ with Haar measure $d z$, set:

$$
A(\Phi)(L, d z)=\int_{L} \Phi(z) d z
$$

And for every $g$ in $G^{1}$, let $d z_{g}$ be the Haar measure on $\rho_{g}(L)$ induced by $d z$. Then the $G^{1}$-stability of $S_{0}\left(F^{2}\right)$ follows from the fact:

$$
A\left(\rho_{g} \Phi\right)(L, d z)=A(\Phi)\left(\rho_{g}(L), d z_{g}\right)
$$

To prove the density of $S_{0}\left(F^{2}\right)$ in $L^{2}\left(F^{2}\right)$, it suffices to show that for every $\Phi$ in $C_{c}^{\infty}\left(F^{2}\right)$, and for any $\varepsilon>0$, we can find a function $\Phi_{\varepsilon}$ in $C_{c}^{\infty}\left(F^{2}\right)$ such that $\Phi-\Phi_{\varepsilon} \in S_{0}\left(F^{2}\right)$ and such that $\left\|\Phi_{\varepsilon}\right\|_{2}<\varepsilon$. Here $\left\|\|_{2}\right.$ is the norm in $L^{2}\left(F^{2}\right)$.

Given $t$ in $F^{*}$, define:

$$
T_{t}(x)=|t|^{-1} \Phi\left(t^{-1} x\right), \text { for all } x \text { in } F^{2}
$$

Then it is straight-forward to see that $A\left(T_{t}(\Phi)\right)=A(\Phi)$ for all $t$ and that $\left\|T_{t}(\Phi)\right\|_{2}=\|\Phi\|_{2}$.

Since $\Phi$ has compact support, we can find a positive number $r$ such that the support of $T_{t}(\Phi)$ is disjoint from the support of $\Phi$ when $|t|>r$. Furthermore, one sees that the functions $\left\{T_{t^{m}}(\Phi) \mid m=0,1,2, \ldots\right\}$ all have disjoint support when $|t|>r$. Fix a $t$ with $|t|>r$, and choose an integer $k$ such that $\sqrt{k} \varepsilon>\|\Phi\|_{2}$. Set:

$$
\Phi_{\varepsilon}=\frac{1}{k} \sum_{j=0}^{k-1} T_{t^{\prime}}(\Phi)
$$

Then certainly $A\left(\Phi_{\varepsilon}\right)=A(\Phi)$. And since $\left\{T_{t^{\prime}}(\Phi)\right\}$ have disjoint support, they are orthogonal in $L^{2}\left(F^{2}\right)$. So we have:

$$
\left\|\Phi_{\epsilon}\right\|_{2}=\frac{1}{k}\left(\sum_{j=1}^{k-1}\left\|T_{t^{\prime}}(\Phi)\right\|_{2}^{2}\right)^{1 / 2}=\frac{1}{k} \sqrt{k\|\Phi\|_{2}^{2}}<\varepsilon,
$$

as desired.
Definition 1.7. A function $h$ on $G^{1}$ is regular $\bmod N$ if
(i) $h$ is $K^{1}$-finite;
(ii) $h(n g)=\eta(n) h(g)$, for all $n$ in $N$ and $g$ in $G^{1}$;
(iii) for every compact set $\Omega$ in $G$, there exist a positive scalar $C$ and a smooth function $\tilde{h}$ on $\Omega$ such that, for all $g$ in $\Omega$, we have:
(A) when $F$ is non-archimedean:

$$
h\left[\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) g\right]= \begin{cases}0, & \text { if }|a|>C^{-1} \\
|a|^{2} \tilde{h}(g), & \text { if }|a|<C\end{cases}
$$

and
(B) when $F$ is archimedean:

$$
h\left[\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) g\right]:\left\{\begin{array}{l}
\rightarrow 0 \text { rapidly as }|a| \rightarrow \infty \\
\left.\langle | a\right|^{\beta_{F}} \cdot \tilde{h}(g), \text { as }|a| \rightarrow 0,
\end{array}\right.
$$

where $\beta_{\mathbf{R}}=2$ and $\beta_{\mathrm{C}}=3 / 2$.
Definition 1.8. $R^{1}(\eta)=\left\{h\right.$ : regular on $\left.G^{1} \bmod N\right\}$.
$R^{1}(\eta)$ is stable under $G$ if $F$ is $p$-adic, and under (Lie $G^{1}, K^{1}$ ) if $F$ is $\mathbf{R}$ or $\mathbf{C}$. (See Wallach [15] for a definition of (Lie $\left.G^{1}, K^{1}\right)$-modules.)

It is simple to check that:

$$
\begin{equation*}
R^{1}(\eta) \subset L^{2}\left(N \backslash G^{1}, \eta\right) . \tag{1.9}
\end{equation*}
$$

Since smooth functions on $G^{1}$ with compact support $\bmod N$ belong to $R^{1}(\eta)$, this space is dense in $L^{2}\left(N \backslash G^{1}, \eta\right)$.

Let $w$ denote the Weyl element $\left(\begin{array}{cc}0 & -1 \\ 1 & 1\end{array}\right)$.
Definition 1.10. For $\Phi$ in $S_{0}\left(F^{2}\right)$, set:

$$
\theta_{\Phi}(g)=\int_{N} f_{\Phi}(w n g) \overline{\eta(n)} d n, \text { for all } g \text { in } G^{1}
$$

It may be checked that this is well-defined. It will turn out that $\theta_{\Phi}$ is even square-integrable $\bmod N$. More precisely, we will show:

Proposition 1.11. $\theta_{\Phi} \in R^{1}(\eta)$, for all $\Phi \in S_{0}\left(F^{2}\right)$.

Proof. Let $\Phi$ be in $S_{0}\left(F^{2}\right)$. Then for all $n$ in $N$ and $g$ in $G$, we have:

$$
\begin{aligned}
\theta_{\Phi}(n g) & =\int_{N} f_{\Phi}(w u n g) \bar{\eta}(u) d u=\int_{N} f_{\Phi}(w u g) \bar{\eta}\left(u n^{-1}\right) d u \\
& =\eta(n) \theta_{\Phi}(g)
\end{aligned}
$$

Observe that, for every $a$ in $F^{*}$, we have:

$$
\theta_{\Phi}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=|a|^{2} \int_{F} \Phi(a, a x) \bar{\psi}\left(a^{2} x\right) d x .
$$

Sending $x$ to $a^{-1} x$, we get:

$$
\theta_{\Phi}\left(\begin{array}{cc}
a & 0  \tag{1.12}\\
0 & a^{-1}
\end{array}\right)=|a| \int_{F} \Phi(a, x) \bar{\psi}(a x) d x
$$

If $F$ is $p$-adic, then $\Phi$ has compact support in $F^{2}$, and hence it is clear that

$$
\theta_{\Phi}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)
$$

vanishes if $|a|$ is large. If $F$ is archimedean, we have for every $M>0$ : $\int_{-M}^{M} \Phi(a, x) \bar{\psi}(a x) d x \rightarrow 0$ rapidly as $|x| \rightarrow \infty$. And we also have: $\Phi(a, x)$ $\rightarrow 0$ rapidly as $|x| \rightarrow \infty$. These two facts imply (using (1.12)) that

$$
\theta_{\Phi}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) \rightarrow 0
$$

rapidly as $|a| \rightarrow \infty$.
Now we look at the behavior when $|a|$ is near 0 . We can write:

$$
\int_{F} \Phi(a, x) \bar{\psi}(a x) d x=f_{1}(a)+f_{2}(a)
$$

where

$$
f_{1}(a)=\int_{F}(\Phi(a, x)-\Phi(0, x)) \bar{\psi}(a x) d x
$$

and

$$
f_{2}(a)=\int_{F} \Phi(0, x)(\bar{\psi}(a x)-1) d x
$$

Assume that $F$ is $p$-adic and that $|a| \ll 1$. Then as $\Phi$ is locally constant, $\Phi(a, x)=\Phi(0, x)$. Consequently $f_{1}(a)=0$. Furthermore, $f_{2}(a)$ equals: $\int_{C} \Phi(0, x)(\bar{\psi}(a x)-1) d x$, where $C$ is a compact set in $F$ outside which $\Phi(0, x)$ is zero. Since $\bar{\psi}$ is a locally constant function of $F$, we have: $\bar{\psi}(a x)=1$ for $x$ in $C$. Hence $f_{2}(a)=\int_{F} \Phi(0, x) d x$, which is zero since $\Phi$
is in $S_{0}\left(F^{2}\right)$. Thus

$$
\theta_{\Phi}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)
$$

vanishes when $|a| \ll 1$.
Let us now turn to the archimedean case. For every $\Phi$ in $S\left(F^{2}\right)$, we set:

$$
\begin{gathered}
\Phi^{\sim}(t)=\int_{F} \Phi(0, x) \bar{\Psi}(t x) d x \\
D_{1} \Phi(u, v)=\frac{\partial \Phi}{\partial u}, \quad \text { and } \quad D \Phi^{\sim}(t)=\frac{d \Phi^{\sim}}{d t}
\end{gathered}
$$

Then we can write: (with $\gamma_{a}$ : a $C^{\infty}$-path from 0 to $a$ in $F$ )

$$
f_{1}(a)=\int_{F} \bar{\psi}(a x) \int_{\gamma_{a}} D_{1} \Phi(t, x) d t=\int_{\gamma_{a}} d t \int_{F} D_{1} \Phi(t, x) \bar{\psi}(a x) d x
$$

and $f_{2}(a)=\int_{\gamma_{a}} D \Phi^{\sim}(t) d t$, where $d t$ is the restriction to $\gamma_{a}$ of the natural 1 -form on $F$. It is clear that both $\int_{F} D_{1} \Phi(t, x) \bar{\psi}(a x) d x$ and $D \phi^{\sim}(t)$ are both $O(1)$ when $|t| \leq|a| \ll 1$. Hence we obtain (for $i=1$ and $i=2$ ):

$$
f_{i}(a)=\left\{\begin{array}{ll}
O(|a|), & \text { if } F=R, \\
O\left(|a|^{1 / 2}\right), & \text { if } F=\mathbf{C},
\end{array} \quad \text { when }|a| \rightarrow 0\right.
$$

It then follows from (1.12) that

$$
\theta_{\Phi}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)
$$

has the desired asymptotic behavior when $|a| \rightarrow 0$.

Remark 1.13. We have actually shown that in the non-archimedean case $\theta_{\Phi}$ is even compactly supported $\bmod N$, for each $\Phi$ in $S_{0}\left(F^{2}\right)$.

## Definition 1.14.

(a) $V_{c}=$ closure of $\left\{\theta_{\Phi} \mid \Phi \in S_{0}\left(F^{2}\right)\right\}$ in $L^{2}(N \backslash G, \eta)$
(b) ${ }^{0} V=$ the orthocomplement of $V_{c}$ in $L^{2}(N \backslash G, \eta)$.

These two spaces are clearly stable under $G^{1}$.
Definition 1.15. For each $h$ in $R^{1}(\eta)$, set:

$$
h^{N}(g)=\int_{N} h\left(w^{-1} n g\right) d n
$$

Proposition 1.16. Let $h$ be in $R^{1}(\eta)$. Then $\int_{N} h\left(w^{-1} n g\right) d n$ is convergent, and the resulting function $h^{N}$ satisfies the following properties:
(a) $h^{N}(n g)=h^{N}(g)$, for all $n$ in $N$ and $g$ in $G$;
(b) $h^{N}$ is $K^{1}$-finite;
(c) (i) when F is p-adic,

$$
h^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) \quad \text { is } 0(\operatorname{resp} . O(1)) \quad \text { if }|a| \ll 1(\operatorname{resp} .|a| \gg 1)
$$

and (ii) when $F$ is archimedean,

$$
h^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) \prec \begin{cases}\left\{\begin{array}{ll}
1, & \text { if } F=\mathbf{R} \\
|a|^{1 / 2}, & \text { if } F=\mathbf{C}
\end{array}\right\}, & \text { as }|a| \rightarrow \infty \\
|a|^{2}, & \text { as }|a| \rightarrow 0\end{cases}
$$

Using this proposition and the identification $N \backslash G^{1} \cong F^{2}-\{0\}$, the following result is easily deduced.

Corollary 1.17. $h^{N} \in L^{2}\left(N \backslash G^{1}\right)$, for all $h \in R^{1}(\eta)$.
Proof of Proposition 1.16. The convergence claim is easy to prove, and the parts (a) and (b) are clear. We give a proof of part (c):
(i) $F$ : p-adic. Let $h$ be in $R^{1}(\eta)$. Then using Iwasawa decomposition, we get:

$$
\begin{align*}
h^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)= & |a|^{2} \int_{\mathscr{D}} h\left[\left(\begin{array}{cc}
a^{-1} & 0 \\
0 & a
\end{array}\right)\left(\begin{array}{cc}
1 & 0 \\
x & 1
\end{array}\right) w^{-1}\right] d x  \tag{1.18}\\
& +|a|^{2} \int_{|x|>1} h\left[\left(\begin{array}{cc}
a^{-1} & 0 \\
0 & a
\end{array}\right)\left(\begin{array}{cc}
1 & -x^{-1} \\
0 & 1
\end{array}\right)\right. \\
& \left.\times\left(\begin{array}{cc}
-x^{-1} & 0 \\
0 & -x
\end{array}\right)\left(\begin{array}{cc}
1 & 0 \\
x^{-1} & 1
\end{array}\right)\right] d x
\end{align*}
$$

Let $\mathrm{I}_{a}$ and $\mathrm{II}_{a}$ respectively denote the first and second expression on the right of (1.18).

When $|a| \ll 1$, since $h$ is regular $\bmod N$,

$$
h\left[\left(\begin{array}{cc}
a^{-1} & 0 \\
0 & a
\end{array}\right) k\right]=0
$$

for all $k$ in $K^{1}$. Hence $\mathrm{I}_{a}=0$ when $|a| \ll 1$.
When $|a| \gg 1$, we have: $\left|a^{-1}\right| \ll 1$ and since $h$ is in $R^{1}(\eta)$ we see that the regularity of $h$ implies that the integrand of $\mathrm{I}_{a}$ is $O\left(|a|^{-1}\right)$. This times $|a|^{2}$ gives $O(1)$. Hence $\mathrm{I}_{a}=O(1)$ when $|a| \gg 1$, since $\operatorname{Vol}(\cap)<\infty$.

It remains to consider $\mathrm{II}_{a}$. Clearly we have

$$
\mathrm{II}_{a}=|a|^{2} \int_{\mathfrak{B}} \psi\left(-a^{2} x\right) h\left[\left(\begin{array}{cc}
-a^{-1} x & 0 \\
0 & -a x^{-1}
\end{array}\right)\left(\begin{array}{cc}
1 & 0 \\
x & 1
\end{array}\right)\right]|x|^{-2} d x
$$

When $|a| \gg 1,\left|a^{-1}\right| \ll 1$ and $\left|a^{-1} x\right|<\left|a^{-1}\right| \ll 1$ for all $x$ in $p$. So the integrand of $\mathrm{II}_{a}$ above is $<\left|a^{-2} x^{2}\right|$, and so:

$$
\mathrm{II}_{a} \prec O(1), \quad \text { when }|a| \gg 1
$$

taking into account the fact that $\operatorname{Vol}(\mathfrak{p})<\infty$.
Now let $|a| \ll 1$. There exist positive scalars $C, D$ and $\delta \in \mathbf{C}$ such that: (for $x$ in $\mathfrak{p}$ )
(A) $\quad \operatorname{supp}:\left\{b \mapsto h\left[\left(\begin{array}{cc}b & 0 \\ 0 & b^{-1}\end{array}\right)\left(\begin{array}{cc}1 & 0 \\ x & 1\end{array}\right)\right]\right\} \subset\{|b| \leq C\}$,
and
(B)

$$
|b| \leq D \Rightarrow h\left(\begin{array}{cc}
b & 0 \\
0 & b^{-1}
\end{array}\right)=\delta|b|^{2}
$$

Consequently, we get:

$$
\mathrm{II}_{a}=|a|^{2} \int_{|x| \leq C|a|} \psi\left(-a^{-1} x\right) h\left[\left(\begin{array}{cc}
-a^{-1} x & 0 \\
0 & a x^{-1}
\end{array}\right)\left(\begin{array}{cc}
1 & 0 \\
x & 1
\end{array}\right)\right]|x|^{-2} d x
$$

Since $|a| \ll 1$ and since $|x| \leq C|a|, h$ is invariant on the right by $\left(\begin{array}{ll}1 & 0 \\ x & 1\end{array}\right)$. Thus we get using (A) and (B):

$$
\begin{aligned}
\mathrm{II}_{a}= & \delta \int_{|x| \leq D|a|} \psi\left(a^{-2} x\right) d x \\
& +|a| \int_{D<|x| \leq C} \psi\left(-x a^{-1}\right) h\left(\begin{array}{cc}
-x & 0 \\
0 & -x^{-1}
\end{array}\right)|x|^{-2} d x
\end{aligned}
$$

Now

$$
\int_{|x| \leq D|a|} \psi\left(-a^{-2} x\right) d x=|a| \int_{|x| \leq D} \psi\left(-a^{-1} x\right) d x=0
$$

if $|a| \ll 1$.
Since $h$ is $K^{1}$-finite, we can write:

$$
h\left(\begin{array}{cc}
x & 0 \\
0 & x^{-1}
\end{array}\right)=\sum_{\omega \in S} h_{\omega}(x)
$$

for all $x$ in $F^{*}$, where $S$ is a finite subset of $\hat{U}_{F}$ and where $h_{\omega}$ satisfies:

$$
h_{\omega}(u x)=\omega(u) h_{\omega}(x), \quad \text { for all } u \text { in } U_{F} \text { and } x \text { in } F^{*}
$$

Consequently we obtain:

$$
\begin{aligned}
& \int_{D<|x| \leq C} \psi\left(-x a^{-1}\right) h\left(\begin{array}{cc}
-x & 0 \\
0 & -x^{-1}
\end{array}\right)|x|^{-2} d x \\
&=\sum_{\left\{n \mid D<q^{-n} \leq C\right\}} q^{2 n} \sum_{\omega \in S} h_{\omega}\left(-\pi^{n}\right) G\left(\omega, a^{-1} \pi^{n}\right)
\end{aligned}
$$

where

$$
\begin{equation*}
G(\omega, t)=\int_{U_{F}} \omega(u) \bar{\psi}(t u) d u, \quad \text { for } t \text { in } F^{*} . \tag{1.19}
\end{equation*}
$$

We will be done if we show that for each $\omega \in S$, and for each $n$ such that $D<q^{-n} \leq C$, the Gauss sum $G\left(\omega, a^{-1} \pi^{n}\right)$ is zero when $|a| \ll 1$. This follows from the following:

Claim 1.20. For every $\omega$ in $\hat{U}_{F}$, there exists $l \in \mathbf{Z}$ such that

$$
|t|>q^{l} \Rightarrow G(\omega, t)=0 .
$$

Indeed, if $\omega=1$ we have: $G(\omega, t)=\int_{Q} \bar{\psi}(t x) d x-\int_{\mathfrak{p}} \bar{\psi}(t x) d x$. And if $\omega \neq 1$ with $\mathfrak{p}^{r}$ the largest ideal such that $\omega$ is trivial on $1+\mathfrak{p}^{r}$, we see that:

$$
G(\omega, t)=\sum_{U_{F} /\left(1+\mathfrak{p}^{\prime}\right)} \omega(u) \bar{\psi}(t u) \int_{\mathfrak{p}^{\prime}} \bar{\psi}(t u x) d x .
$$

If $p^{-n(\omega)}=\operatorname{ker}(\psi)$, we get by orthogonality of characters of compact groups that:

$$
G(\omega, t)=0, \quad \text { if } v(t)<-l_{\omega}
$$

where

$$
l_{\omega}=n(\psi)+1(\text { resp. } n(\psi)+r) \quad \text { if } \omega=1(\text { resp. } \omega \neq 1) .
$$

(ii) $F$ : archimedean. We will first study the behavior of

$$
h^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) \quad \text { as }|a| \rightarrow 0 .
$$

Since $h$ is $K^{1}$-finite, there exists an $f$ in $C_{c}^{\infty}\left(G^{1}\right)$ such that $h=h * f^{\vee}$. Here $f^{\vee}$ denotes the function: $g \mapsto f\left(g^{-1}\right)$, and $*$ denotes convolution product. It then follows that:

$$
\begin{align*}
h^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) & =\left(h^{N} * f^{\vee}\right)\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)  \tag{1.21}\\
& =\int_{M G^{1}} h^{N}\left[\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) g\right] \mu_{f}(g) d g,
\end{align*}
$$

where $\mu_{f}(g)=\int_{N \backslash G^{1}} f(n g) d n$.

Note that $\mu_{f}$ has compact support $\bmod N$. So there exists a compact set $\Omega_{1}$ in $A^{1} K^{1} \simeq N \backslash G^{1}$ such that:

$$
h^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=\iint_{\Omega_{1}} h^{N}\left[\left(\begin{array}{cc}
a b & 0 \\
0 & a^{-1} b^{-1}
\end{array}\right) d\right] \cdot \mu_{f}\left[\left(\begin{array}{cc}
b & 0 \\
0 & b^{-1}
\end{array}\right) k\right] \frac{d^{*} b d k}{|b|^{2}} .
$$

Since $\mu_{f} \prec 1$, this yields:

$$
\left|h^{N}\left(\begin{array}{cc}
a & 0  \tag{1.22}\\
0 & a^{-1}
\end{array}\right)\right|<\iint_{\Omega_{1}}\left|h^{N}\left[\left(\begin{array}{cc}
a b & 0 \\
0 & a^{-1} b^{-1}
\end{array}\right) k\right]\right| \cdot|b|^{-1} d^{*} b d k
$$

Let $\Omega_{2}(a)$ be the image of $\Omega_{1}$ under the map:

$$
\left(\begin{array}{cc}
b & 0 \\
0 & b^{-1}
\end{array}\right) k \mapsto\left(\begin{array}{cc}
a^{-1} b & 0 \\
0 & a b^{-1}
\end{array}\right) k
$$

Then (1.22) becomes:

$$
\left|h^{N}\left(\begin{array}{cc}
a & 0  \tag{1.23}\\
0 & a^{-1}
\end{array}\right)\right|<|a|^{2} \cdot H(a),
$$

where

$$
H(a)=\iint_{\Omega_{2}(a)}\left|h^{N}\left(\begin{array}{cc}
b & 0 \\
0 & b^{-1}
\end{array}\right) k\right| \cdot|b|^{-1} d^{*} b d k .
$$

Observing that $H(a)$ is at least $O(1)$ for $|a|$ small, we get from (1.23):

$$
\left|h^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)\right|<|a|^{2} \quad \text { as }|a| \rightarrow 0 .
$$

It remains to show the asserted behavior of

$$
h^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)
$$

as $|a| \rightarrow \infty$ in the archimedean situation. For this we use the Iwasawa decomposition: $G^{1}=N A^{1} K^{1}$ and write:

$$
\left(\begin{array}{ll}
1 & 0 \\
x & 1
\end{array}\right)=\left(\begin{array}{cc}
1 & -\bar{x}(1-x \bar{x})^{1 / 2} \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
(1-x \bar{x})^{-1 / 2} & 0 \\
0 & (1-x \bar{x})^{1 / 2}
\end{array}\right) k_{x},
$$

where

$$
k_{x}=\left(\begin{array}{cc}
\bar{x}(1-x \bar{x})^{-1 / 2} & -(1-x \bar{x})^{-1 / 2} \\
(1-x \bar{x})^{-1 / 2} & x(1-x \bar{x})^{-1 / 2}
\end{array}\right) \in K^{1} .
$$

(Note that when $F=\mathbf{R}, x=\bar{x}$.)

Thus:

$$
\left.\begin{array}{rl}
h^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=|a|^{2} \int_{F} h\left[\left(\begin{array}{cc}
a^{-1}(1-x \bar{x})^{-1 / 2} & 0 \\
0 & a(1-x \bar{x})^{1 / 2}
\end{array}\right) k_{x} w^{-1}\right.
\end{array}\right] .
$$

Again, when $|a| \rightarrow \infty,|a|^{-1} \cdot\left|(1-x \bar{x})^{-1 / 2}\right| \rightarrow 0$, and the regularity of $h$ implies:

$$
h\left[\left(\begin{array}{cc}
a^{-1}(1-x \bar{x})^{-1 / 2} & 0 \\
0 & a(1-x \bar{x})^{1 / 2}
\end{array}\right) k_{x} w^{-1}\right] \prec|a|^{-\beta_{F}} \cdot(1-x \bar{x})^{-\beta_{F} / 2}
$$

Consequently,

$$
\begin{aligned}
h^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) \prec|a|^{2-\beta_{F}} \int_{F}(1-x \bar{x})^{-\beta_{F} / 2} d x & \prec|a|^{2-\beta_{F}}, \\
\text { when }|a| & \rightarrow \infty, \text { as desired. }
\end{aligned}
$$

Proposition 1.24. If $h$ is in $R^{1}(\eta)$, then

$$
h \in{ }^{0} V \Leftrightarrow h^{N}=0 .
$$

Proof. For every $\Phi$ in $S_{0}\left(F^{2}\right)$ and every $h$ in $R^{1}(\eta)$, a formal manipulation yields:

$$
\begin{equation*}
\left\langle\theta_{\Phi}, h\right\rangle=\left\langle f_{\Phi}, h^{N}\right\rangle \tag{1.25}
\end{equation*}
$$

To make this rigorous we need only to check that: $\int_{G^{1}} f(w g) \overline{h(g)} d g$ is absolutely convergent. We may replace $h$ by $|h|$ and assume that $h \geq 0$. Then the desired result follows from the asymptotic properties of $f_{\Phi}$ and $h$ (cf. Proposition 1.16).

Since $S_{0}\left(F^{2}\right)$ is dense in $L^{2}\left(F^{2}\right)$, we see that $S_{0}\left(N \backslash G^{1}\right)=\left\{f_{\Phi} \mid \Phi \in\right.$ $\left.S_{0}\left(F^{2}\right)\right\}$ is dense in $L^{2}\left(N \backslash G^{1}\right)$. Therefore any $h$ in $R^{1}(\eta)$ is orthogonal to $V_{c}=$ the closure of $\left\{\theta_{\Phi} \mid \Phi \in S_{0}\left(F^{2}\right)\right\}$ if and only if $h^{N}$ is identically zero.

Definition 1.26.

$$
{ }^{0} V_{K}=\left\{h \in{ }^{0} V \mid h \text { is } K^{1} \text {-finite }\right\} .
$$

Theorem 1.27.
(a) ${ }^{0} V_{K}$ is admissible;
(b) ${ }^{0} V_{K}={ }^{0} V \cap R^{1}(\eta)$, i.e., every $h$ in ${ }^{0} V_{K}$ is regular $\bmod N$; and
(c) ${ }^{0} V_{K}$ is precisely the space of all $K^{1}$-finite vectors in the discrete spectrum of $L^{2}\left(N \backslash G^{1} ; \eta\right)$.

Remark 1.28. It is known (cf. [10], for example) that $L^{2}(N \backslash G ; \eta)$ is a multiplicity-free representation. Hence every irreducible of $G^{1}$ occurring in ${ }^{0} V_{K}$ will occur exactly once.

Proof of Theorem 1.27. We will show the admissbility (a) and regularity (b) of ${ }^{0} V_{K}$ simultaneously.

We first decompose ${ }^{0} V_{K}$ according to the right action as:

$$
{ }^{0} V_{K}=\underset{\left(\sigma, H_{\sigma}\right) \in \widehat{K^{1}}}{ }{ }^{0} V_{K}(\sigma),
$$

where each $\sigma$-isotypic space ${ }^{0} V_{K}(\sigma)$ can be identified with

$$
\begin{aligned}
& \left\{f: G^{1} \rightarrow H_{\sigma} \mid f(n g k)=\eta(n) \sigma\left(k^{-1}\right) f(g), \text { for all } n \in N,\right. \\
& \\
& \left.\qquad g \in G^{1} \text { and } k \in K^{1} ; \text { and } f \in L^{2}\left(N \backslash G^{1} / K\right)\right\} .
\end{aligned}
$$

For every $\left(\sigma, H_{\sigma}\right)$ in $\widehat{K^{1}}$ and for each character $\omega$ of $B_{K}^{1}=B \cap K^{1}$, we let $H_{\sigma}(\omega)$ denote the subspace of $H_{\sigma}(\omega)$ consisting of vectors $v$ such that $\sigma(b) v=\omega(b) v$ for all $b$ in $B_{K}^{1}$. Then it can be seen (cf. [5]) that $\operatorname{dim} H_{\sigma}(\omega) \leq 1$, and that, for a given $\sigma, H_{\sigma}(\omega)=0$ for all but finitely many characters $\omega$ of $B_{K}^{1}$. We get:

$$
{ }^{0} V_{K}(\sigma)=\bigoplus_{\omega \in B_{K}^{1}}{ }^{0} V_{K}(\sigma, \omega)
$$

where

$$
{ }^{0} V_{K}(\sigma, \omega)=\left\{f \in{ }^{0} V_{K}(\sigma) \mid f(g b)=\bar{\omega}(b) f(g)\right.
$$

for all $g$ in $G^{1}$ and $b$ in $\left.B_{K}^{1}\right\}$.
To show the admissibility of ${ }^{0} V_{K}$, we need to prove that ${ }^{0} V_{K}(\sigma, \omega)$ is finite-dimensional for every $\left(\sigma, H_{\sigma}\right) \in \widehat{K^{1}}$ and for every $\omega$ is $\widehat{B_{K}^{1}}$.
$F$ : non-archimedean
The desired result in this case will follow immediately from:
Lemma 1.29. Fix $\left(\sigma, H_{\sigma}\right) \in \hat{K}_{1}$ and $\omega \in \hat{B}_{k}^{1}$. Then we can find $a$ positive constant $C=C(\sigma, \omega)$ such that for every $h$ in ${ }^{0} V_{K}(\sigma, \omega)$ there exists a vector $v$ in $H_{\sigma}(\omega)$ with the property:

$$
h\left[\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)\right]= \begin{cases}|a|^{2} \bar{\omega}\left(\begin{array}{cc}
a /|a| & 0 \\
0 & |a| / a
\end{array}\right) v_{0} & \text { if }|a| \leq C \\
0 & \text { if }|a| \geq C^{-1}\end{cases}
$$

For every integer $r \geq 1$, let $K_{r}^{1}$ denote the principal congruence subgroup of level $r$, namely: $\left\{k \in K^{1} \mid d \mathrm{k} \equiv I\left(\bmod \mathfrak{p}^{r}\right)\right\}$. Let $m$ be the smallest integer $r$ such that $K_{r}^{1} \subset \operatorname{Ker}(\sigma)$.

Let $x$ be in $\mathfrak{p}^{m}$. Then, for every $h$ in ${ }^{0} V(\sigma, \omega)$, we have:

$$
h\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=h\left[\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)\left(\begin{array}{cc}
1 & x \\
0 & 1
\end{array}\right)\right], \quad \text { for all } a \text { in } F^{*}
$$

But

$$
\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)\left(\begin{array}{cc}
1 & x \\
0 & 1
\end{array}\right)=\left(\begin{array}{cc}
1 & a^{2} x \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)
$$

so

$$
h\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=\psi\left(a^{2} x\right) h\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right), \quad \forall a \in F^{*}, x \in \mathfrak{p}^{m}
$$

Let $\mathfrak{p}^{-n(\psi)}$ be the largest ideal on which $\psi$ is trivial. Then we get:

$$
h\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=0, \quad \text { if } 2 v(a)<-m-n(\psi)
$$

i.e., if $|a|>q^{m+n(\psi) / 2}$.

Since the case when $H_{\sigma}(\omega)=0$ is trivial, we will choose $\left(\sigma, H_{\sigma}\right)$ and $\omega$ such that $H_{\sigma}(\omega) \neq\{0\}$. We will produce certain explicit special functions $\theta_{\Phi}$ in $V_{c}(\sigma, \omega)$ so that $h$ is forced to have the derived asymptotic property as $|a| \ll 1$ by virtue of $h$ being orthogonal to $V_{c}(\sigma, \omega)$. For any right $G^{1}$-module $X$, the space $X(\sigma, \omega)$ is defined in the obvious way.

For any $\Phi$ in $S_{0}\left(F^{2} ; \sigma, \omega\right)$, we have:

$$
\begin{aligned}
\theta_{\Phi}\left(\begin{array}{cc}
\pi^{l} & 0 \\
0 & \pi^{-l}
\end{array}\right)= & q^{-l} \int_{F} \Phi\left(\pi^{l}, x\right) \bar{\psi}\left(\pi^{l} x\right) d x \\
= & q^{-l} \int_{p^{\prime}} \Phi\left[\left(0, \pi^{l}\right)\left(\begin{array}{cc}
0 & -1 \\
1 & \pi^{-l} x
\end{array}\right)\right] \bar{\psi}\left(\pi^{l} x\right) d x \\
& +q^{-l} \int_{|x|>q^{-1}} \Phi\left[(0, x)\left(\begin{array}{cc}
1 & 0 \\
\pi^{l} x^{-1} & 1
\end{array}\right)\right] \bar{\psi}\left(\pi^{l} x\right) d x
\end{aligned}
$$

We will choose $\Phi$ in such a way that the support of $\{t \mapsto \Phi[(0, t) k]\}$ is in $\left\{t^{-1} \in \mathfrak{p}^{m-l}\right\}$, uniformly for all $k$ in $K^{1}$. Then the integral:

$$
\int_{\mathfrak{p}^{\prime}} \Phi\left[\left(0, \pi^{l}\right)\left(\begin{array}{cc}
0 & -1 \\
1 & \pi^{-l} x
\end{array}\right)\right] \bar{\psi}\left(\pi^{l} x\right) d x
$$

vanishes. Furthermore if $x^{-1} \in \mathfrak{p}^{m-l}$, then the matrix

$$
\left(\begin{array}{cc}
1 & 0 \\
\pi^{\prime} x^{-1} & 1
\end{array}\right)
$$

is in $K_{m}^{1}$, and so leaves $\Phi$ invariant (on the right). Thus we get:

$$
\theta_{\Phi}\left(\begin{array}{cc}
\pi^{l} & 0  \tag{1.30}\\
0 & \pi^{-l}
\end{array}\right)=q^{-1} \int_{|x|>q^{m-1}} \Phi(0, x) \bar{\psi}\left(\pi^{l} x\right) d x=q^{-l} \cdot \Phi^{\sim}\left(\pi^{l}\right)
$$ where $\Phi^{-}(u)=\int_{F} \Phi(0, x) \bar{\psi}(u x) d x$.

Let $\omega \neq 1$. Then the conductor of $\omega$ is $\mathfrak{p}^{r}$ with $r>0$, i.e., $1+\mathfrak{p}^{r}$ is the largest subgroup of $U_{F}$ on which $\omega$ is trivial. Choose a non-zero vector $v_{0}$ in $H_{\sigma}(\omega)$ and define for every $j \geq[3 m / 2]-n(\psi)$ :

$$
\phi_{j}\left(\pi^{-i}\right)=\delta_{i, j+r+n(\psi)} \cdot v_{0} .
$$

Since every non-zero element of $F^{2}-\{0\}$ can be written as $\left(0, \pi^{-i}\right) k$ for some integer $i$ and some $k$ in $K^{1}$, we see that we get a unique function $\Phi_{j}$ in $S_{0}\left(F^{2} ; \sigma, \omega\right)$ such that $\Phi_{j}\left(0, \pi^{-i}\right)$. Then for every integer $l \geq[-m / 2]$, we have:

$$
l+j+r+n(\psi)>l+j+n(\psi) \geq m-1
$$

and thus:

$$
\operatorname{Supp}\left(\phi_{j}\right) \subset\left\{\pi^{l} x^{-1} \in \mathfrak{p}^{m}\right\}
$$

So by (1.30), we get:

$$
\begin{aligned}
\theta_{\Phi_{J}}\left(\begin{array}{cc}
\pi^{l} & 0 \\
0 & \pi^{-l}
\end{array}\right) & =q^{-l+j+r+n(\psi)} \cdot v_{0} \int_{U_{F}} \bar{\psi}\left(\pi^{l-j-r-n(\psi)} \cdot u\right) \omega(u) d u \\
& =\delta_{l, j} \cdot q^{r} \cdot \operatorname{Vol}\left(U_{F}\right)
\end{aligned}
$$

Now, since $h$ is orthogonal to $\theta_{\Phi_{j}}$, we have for every $j \geq[3 m / 2]$ $n(\psi):$

$$
\begin{aligned}
0 & =\int_{|a| \leq q^{m / 2}}\left\langle h\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right), \theta_{\Phi_{j}}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)\right)|a|^{-2} d^{*} a \\
& =\sum_{l \geq[-m / 2]} q^{2 l}\left(h\left(\begin{array}{cc}
\pi^{l} & 0 \\
0 & \pi^{-l}
\end{array}\right), \theta_{\Phi_{j}}\left(\begin{array}{cc}
\pi^{l} & 0 \\
0 & \pi^{-l}
\end{array}\right)\right\rangle .
\end{aligned}
$$

Consequently we get:

$$
q^{2 j+r}\left(\phi\left(\begin{array}{cc}
\pi^{j} & 0 \\
0 & \pi^{-j}
\end{array}\right), V_{0}\right) \cdot \operatorname{Vol}\left(U_{F}\right)=0
$$

$$
\text { for every } j \geq[3 m / 2]-n(\psi)
$$

Hence: $h\left(\begin{array}{cc}a & 0 \\ 0 & a^{-1}\end{array}\right)=0, \quad$ if $|a| \leq q^{n(\psi)-3 m / 2}$.

We now turn to the case when $\omega=1$.
Choose for every $j>3 m / 2-n(\psi)$, a function $\Phi_{j} \in S_{0}\left(F^{2} ; \sigma, 1\right)$ such that:

$$
\Phi_{j}^{\sim}\left(\pi^{l}\right)= \begin{cases}q^{r+n(\psi)} \cdot v_{0}, & \text { if } l=j \\ -q^{r+n(\psi)-1} \cdot v_{0}, & \text { if } l=j+1 \\ 0, & \text { otherwise }\end{cases}
$$

Once again $v_{0}$ is some fixed non-zero vector in $H_{\sigma}(1)$.

$$
\begin{aligned}
& \left\langle h, \theta_{\Phi_{J}}\right\rangle=0 \text { implies: } \\
& \left\langle h\left(\begin{array}{cc}
\pi^{j} & 0 \\
0 & \pi^{-1}
\end{array}\right), v_{0}\right\rangle q^{j}=\left\langle h\left(\begin{array}{cc}
\pi^{j+1} & 0 \\
0 & \pi^{-j-1}
\end{array}\right), v_{0}\right\rangle \cdot q^{j+2} \\
& \text { for every } j>3 m / 2-n(\psi) \text {. }
\end{aligned}
$$

Then we must have:

$$
h\left(\begin{array}{cc}
\pi^{j} & 0 \\
0 & \pi^{-j}
\end{array}\right)=-q^{-2 j} \cdot v_{0}, \quad \forall j>\frac{3 m}{2}-n(\psi)
$$

Since $h$ is invariant under the right action of $B_{K}^{1}$, we get:

$$
h\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=|a|^{2} v_{0}, \quad \text { if }|a| \leq q^{n(\psi)-3 m / 2}
$$

This proves Lemma 1.29 and hence parts (a) and (b) of Theorem 1.27 in the non-archimedean case.
$F=\mathbf{R}$
For each integer $n$, let $\sigma_{n}$ denote the irreducible module of $\mathrm{SO}(2)$ given by: $k_{\theta} \mapsto e^{i n \theta}$, where

$$
k_{\theta}=\left(\begin{array}{rr}
\cos \theta & -\sin \theta \\
\sin \theta & \cos \theta
\end{array}\right)
$$

It suffices to show that each ${ }^{0} V_{K}\left(\sigma_{n}\right)$ is finite-dimensional. This will follow from the following explicit bound:

Lemma 1.31.

$$
\operatorname{dim}_{\mathbf{C}}^{0} V_{K}\left(\sigma_{n}\right) \leq \begin{cases}\max \left\{0, \frac{|n|-3}{2}\right\}, & \text { if } n \equiv 1(2) \\ \frac{|n|}{2}, & \text { if } n \equiv 0(2)\end{cases}
$$

Proof. Choose for every $l \in \mathbf{Z}^{+}$, a function $\Phi_{l}$ in $S\left(F^{2} ; \sigma_{n}\right)$ by:

$$
\Phi_{l}(0, x)=P_{l} e^{-\pi x^{2}}, \quad \text { where } P_{l}=\sum_{j=0}^{l} c_{j l} x^{2 j+m}, m=|n|
$$

such that $c_{l /} \neq 0$. Since $\mathbf{R}^{2}-\{0\}=Y \cdot K^{1}$ with $Y=\left\{(0, x) \mid x \in \mathbf{R}^{*}\right\}$, every $\Phi$ in $S\left(F^{2} ; \sigma_{n}\right)$ is determined by its restriction to $Y$.

For $\Phi$, to be in $S_{0}\left(\sigma_{n}\right)$, we need:

$$
\sum_{j=0}^{l} c_{j l} \int_{\mathbf{R}} x^{2 \jmath+m} e^{-\pi x^{2}} d x=0
$$

which is automatic if $m$ is odd. When $m$ is even, choose $\left\{c_{j l}\right\}$ such that

$$
\sum_{0 \leq j \leq l} c_{j l} \pi^{-j} \int_{\mathbf{R}_{+}^{*}} x^{j+(m+1) / 2} e^{-x} d^{*} x=0
$$

i.e., $\Sigma_{0 \leq j \leq l} c_{j l} \pi^{-j} \Gamma(j+(m+1) / 2)=0$. We can do this for every $l \geq 0$. Set $\boldsymbol{\theta}_{l}=\boldsymbol{\theta}_{\Phi_{l}}$, and let $a \in \mathbf{R}^{*}$. Then we obtain:

$$
\theta_{l}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=|a|^{2} \int_{\mathbf{R}} \Phi_{l}\left[\left(0, a\left(1+x^{2}\right)^{1 / 2}\right) k_{x}\right] \bar{\psi}\left(a^{2} x\right) d x
$$

where

$$
k_{x}=\left(\begin{array}{cc}
x\left(1+x^{2}\right)^{-1 / 2} & -\left(1-x^{2}\right)^{-1 / 2} \\
\left(1+x^{2}\right)^{-1 / 2} & x\left(1+x^{2}\right)^{-1 / 2}
\end{array}\right) \in \mathrm{SO}(2)
$$

Thus $\theta_{l}\left(\begin{array}{cc}a & 0 \\ 0 & a^{-1}\end{array}\right)$ equals:

$$
\begin{array}{r}
|a|^{2} \sum_{j=0}^{l} c_{j l} \int_{\mathbf{R}} a^{2 j+m}\left(1+x^{2}\right)^{j+m / 2} e^{-\pi a^{2}\left(1+x^{2}\right)} \frac{(x \pm i)^{m}}{\left(1+x^{2}\right)^{m / 2}} \bar{\psi}\left(a^{2} x\right) d x \\
\text { where } \pm \text { denotes } \operatorname{sgn}(n), \text { which becomes: } \\
|a| e^{-\pi a^{2}} \sum_{0 \leq j \leq l} c_{l l} \int_{\mathbf{R}}(x \pm i a)^{m+j}(x \mp i a)^{j} e^{-\pi x^{2}} \bar{\psi}(a x) d x \\
\text { as } x \rightarrow a^{-1} x
\end{array}
$$

Let
$Q_{l}(a)=\sum_{0 \leq j \leq l} c_{l l} \sum_{p=0} \sum_{q=0}\binom{j}{p}\binom{m+j}{j}( \pm i)^{q-p-m} a^{m+2 j-p-q} H_{p+q}(\sqrt{2 \pi a})$,
where $H_{p+q}$ is the $(p+q)$ th Hermite polynomial (see Magnus [9]). Then we obtain:

$$
\theta_{l}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=|a| e^{-2 \pi a^{2}} Q_{l}(a)
$$

Let $d_{l}$ be the coefficient of the leading term $a^{m+2 l}$ in $Q_{l}$. Then

$$
\begin{aligned}
d_{l} & =c_{j l} \sum_{p=0}^{l} \sum_{q=0}^{m+l}\binom{l}{p}\binom{m+l}{q}( \pm i)^{q-p-m_{2} p+q}(\sqrt{2 \pi})^{p+q} \\
& =c_{l l}( \pm i)^{m} \sum_{k=0}^{m+2 l} \beta_{k l}( \pm \sqrt{-8 \pi})^{k},
\end{aligned}
$$

where

$$
\beta_{k l}=\sum_{p=0}^{\min (l, k)}(-i)^{p}\binom{l}{p}\binom{m+l}{k-p} .
$$

Clearly, for every $l$, there exists a $k$ such that $\beta_{k l} \neq 0$. And $c_{l l} \neq 0$ by assumption. Thus $d_{l} \neq 0$, for otherwise $\sqrt{\pi}$ would be algebraic! Therefore, the degree $Q_{l}$ is precisely $m+2 l$.

Let $\phi \in{ }^{0} V_{K}\left(\sigma_{n}\right)$. Then $\phi \perp\left\{\theta_{l}\right\}$, and we have

$$
\begin{aligned}
0 & =\int_{\mathbf{R}^{*}} \phi\left(\begin{array}{cc}
0 & 0 \\
0 & a^{-1}
\end{array}\right) \overline{Q_{l}(a)}|a| e^{-2 \pi a^{2}}|a|^{-2} d^{x} a \\
& =\int_{\mathbf{R}} \tilde{\phi}(a) Q_{l}^{\prime}(a) e^{-2 \pi a^{2}} d a,
\end{aligned}
$$

where

$$
\tilde{\phi}(a)= \begin{cases}\phi\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right), & \text { if } a \neq 0 \\
0 & \text { if } a=0\end{cases}
$$

and $Q_{l}^{\prime}(a)=a^{-2} \overline{Q_{l}(a)}$. Note that $Q_{l}^{\prime}$ is a polynomial of degree $m+2 l-$ 2 if $m+2 l \geq 2$.

Let $Z_{m}$ be the closure in $L^{2}(\mathbf{R})$ of the span of the union of $\left\{a^{k} e^{-\pi a^{2}}\right\}$ $\left.k \in \mathbf{Z}_{+}, k-m \equiv 1(2)\right\}$ and

$$
\left\{Q_{l}^{\prime}(a) e^{-2 \pi a^{2}} \left\lvert\, \begin{array}{l}
l>0 \text { if } m: \text { even or } m=1, \\
l \geq 0 \text { if } m: \text { odd and } \geq 3
\end{array}\right.\right\} .
$$

Let $Z_{m}^{\perp}$ denote the ortho-complement of $Z_{m}$ in $L^{2}(\mathbf{R})$.
Then

$$
\operatorname{dim}_{\mathbf{C}} Z_{m}^{\perp} \leq \begin{cases}\max \{0, m-3 / 2\} & \text { if } m \equiv 1(2) \\ \frac{m}{2} & \text { if } m \equiv 0(2) .\end{cases}
$$

Claim. $\tilde{\phi} \in L^{2}(R)$, for all $\phi \in^{0} V\left(\sigma_{n}\right)$.

We are done modulo this claim because then $\left\{\tilde{\phi} \mid \phi \in{ }^{0} V\left(\sigma_{n}\right)\right\} \subset Z_{m}^{\perp}$ and $\phi \rightarrow \tilde{\phi}$ is injective.

To prove the claim, observe that since $\phi$ is $K$-finite, there exists $f \in C_{c}^{\infty}(G)$ such that $\phi=\phi * f^{\vee}$. Hence

$$
\begin{aligned}
\tilde{\phi}(a)= & \left(\phi * f^{\vee}\right)\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=\int_{G} \phi\left(\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) g\right) f(g) d g \\
= & \int_{N \backslash G} \phi\left(\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) g\right) d g \\
& \cdot \int_{N} \psi\left(\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) n\left(\begin{array}{cc}
a^{-1} & 0 \\
0 & a
\end{array}\right)\right) f(n g) d n .
\end{aligned}
$$

There exists a compact set $\Omega$ in $N \backslash G \simeq A K$ such that the inner integral vanishes unless $g \in \Omega$. So

$$
\tilde{\phi}(a)=\int_{\Omega} \phi\left(\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) g\right) \hat{f}_{g}\left(a^{2}\right) d g
$$

where

$$
\hat{f}_{g}(y)=\int_{\mathbf{R}} f\left(\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right) g\right) \psi(x y) d x
$$

There exists a positive function $h$ in $S(\mathbf{R})$ such that $\left|\hat{f}_{g}\right| \leq h$, for $g \in \Omega$. This implies

$$
\begin{aligned}
|\phi(a)| & \leq h\left(a^{2}\right) \int_{\Omega}\left|\phi\left(\left(\begin{array}{cc}
a b & 0 \\
0 & a^{-1} b^{-1}
\end{array}\right) k\right)\right| \cdot|b|^{-2} d^{*} b d k \\
& \prec h\left(a^{2}\right) \int_{\omega}\left|\phi\left(\begin{array}{cc}
a b & 0 \\
0 & a^{-1} b^{-1}
\end{array}\right)\right||b|^{-2} d^{*} b
\end{aligned}
$$

where $\omega$ is a compact set in $A$

$$
=h\left(a^{2}\right) \cdot|a|^{2} \int_{\omega}\left|\left(\begin{array}{cc}
b & 0 \\
0 & b^{-1}
\end{array}\right)\right| \cdot|b|^{-2} d^{*} b \prec h(a) \cdot|a|^{2} .
$$

This proves the claim.
Incidentally we have also shown the regularity $\bmod N$ of functions in ${ }^{0} V\left(\sigma_{n}\right)$.

$$
F=\mathbf{C}
$$

In this case the unitary irreducibles of $K^{1}=\mathrm{SU}(2)$ are given by:

$$
\sigma_{m}: K^{1} \rightarrow \text { Aut } H_{m}, \quad m \in \mathbf{Z}_{+},
$$

where $H_{m}=$ \{homogeneous polynomials of degree $m$ in $\left.\mathbf{C}[u, v]\right\}$, with $u$, $v$ : indeterminates over $\mathbf{C}$.

The characters $\omega_{n}$ of

$$
B_{K}^{1}=\left\{\left(\begin{array}{cc}
e^{i \theta} & 0 \\
0 & e^{-i \theta}
\end{array}\right)\right\}
$$

are given by:

$$
\left(\begin{array}{cc}
e^{i \theta} & 0 \\
0 & e^{-i \theta}
\end{array}\right) \leftrightarrow e^{i n \theta}, \quad \text { for } n \in \mathbf{Z} .
$$

It can be seen that $H_{m}\left(\omega_{n}\right) \neq 0$ iff $|n| \leq m$ and $m \equiv|n|(\bmod 2)$.
Explicitly, $H_{m}\left(\omega_{n}\right)$ is the $\mathbf{C}$-span of $u^{m_{1}} \cdot v^{m_{2}}$, where

$$
m_{1}=\frac{m+|n|}{2} \quad \text { and } \quad m_{2}=\frac{m-|n|}{2} .
$$

We will be done if we show:
Lemma 1.32.
(a) $\operatorname{dim}_{\mathbf{C}}{ }^{0} V_{K}\left(\sigma_{m}, \omega_{n}\right) \leq m+1$, for all $m \in \mathbf{Z}_{+}, n \in \mathbf{Z}$.
(b) The functions in ${ }^{0} V_{K}\left(\sigma_{m}\right)$ are regular on $G^{1} \bmod N$.

The proof of this lemma is very similar to that of Lemma 1.31 (the real case), and will be omitted.

It remains to prove part (c) of Theorem 1.27.
Definition 1.33. $V_{\text {disc }}=\{K$-finite vectors in the discrete spectrum of $\left.L^{2}\left(N \backslash G^{1}, \eta\right)\right\}$.

Since we have shown that ${ }^{0} V_{K}$ is an admissible subspace of $L^{2}\left(N \backslash G^{1}, \eta\right)$, we have:

$$
\begin{equation*}
{ }^{0} V_{K} \subset V_{\text {disc }} . \tag{1.34}
\end{equation*}
$$

In what follows we will show the reverse inclusion.
Let ( $\pi, H_{\pi}$ ) be an irreducible unitary representation of $G^{1}$. Let $H_{\pi, K}$ denote the space of $K$-finite vectors in $H_{\pi}$. Then we get an admissible representation of $G^{1}$ (resp. (Lie $G^{1}, K^{1}$ )) on $H_{\pi, K}$ when $F$ is $p$-adic (resp. archimedean).

Definition 1.35. An irreducible, square-integrable $G^{1}$-module ( $\pi, H_{\pi}$ ) has a Whittaker model relative to $\eta$ iff $\operatorname{Hom}_{X}\left(H_{\pi, K}, V_{\text {disc }}\right) \neq\{0\}$, where

$$
X= \begin{cases}G^{1}, & \text { if } F \text { is } p \text {-adic } \\ \left(\operatorname{Lie} G^{1}, K^{1}\right), & \text { if } F \text { is archimedean. }\end{cases}
$$

For a definition of intertwining operators between (Lie $G^{1}, K^{1}$ )-modules, see Wallach [[15]).

## Definition 1.36.

$$
\sum(\eta)=\left\{\begin{array}{r}
\text { Irreducible, inequivalent, square-integrable } G^{1} \text {-modules } \\
\text { which have Whittaker models relative to } \eta
\end{array}\right\}
$$

Remark 1.37. It is known that the right regular representation $L^{2}(G)$ is a multiple of (the multiplicity-free representation) $L^{2}(N \backslash G, \eta)$. This is a theorem of B. Blackadar, a proof of which can be found in [6]. From this theorem it can be deduced that $V_{\text {disc }}$ embeds in the discrete spectrum of $L^{2}\left(G^{1}\right)$. Thus we have the following decomposition:

$$
V_{\mathrm{disc}} \simeq \bigoplus_{\pi \in \Sigma(\eta)} \pi
$$

Let $\left(\pi, H_{\pi}\right)$ be in $\Sigma(\eta)$. Then, since $V_{\text {disc }}$ is multiplicity-free, $H_{\pi, K}$ embeds in $V_{\text {disc }}$ in a unique manner.

Definition 1.38. Let $\left(\pi, H_{\pi}\right)$ be in $\Sigma(\eta)$.
(a) $\mathfrak{W}(\pi)=$ the unique component of $V_{\text {disc }}$ isomorphic to $H_{\pi, K}$.
(b) A Whittaker vector (relative to $\eta$ ) associated to ( $\pi, H_{\pi}$ ) is an element of $\mathfrak{W}(\pi)$.

Since we have the decomposition:

$$
\begin{equation*}
V_{\mathrm{disc}}=\bigoplus_{\left(\pi, H_{\pi}\right) \in \Sigma(\eta)} \mathfrak{W}(\pi), \tag{1.39}
\end{equation*}
$$

we need only to show that $\mathfrak{W}(\pi)$ is in ${ }^{0} V_{K}$ for all $\left(\pi, H_{\pi}\right)$ in $\Sigma(\eta)$. This will be achieved by the following:

Proposition 1.40. Let $\left(\pi, H_{\pi}\right)$ be in $\Sigma(\eta)$. Then every Whittaker vector $W$ (relative to $\eta$ ) associated to $\pi$ is regular on $G^{1} \bmod N$, and satisfies: $W^{N} \equiv 0$.

## Proof.

F: p-adic
Every discrete series representation of $G^{1}$ is either supercuspidal or special (cf. [4], [7] for definition). The proposition is a standard fact (cf. [7]) for any supercuspidal $\pi$ in $\Sigma(\eta)$. Assume that $\pi$ is special. Then we know from Godement ([4]) that every Whittaker vector $W$ associated to $\pi$ satisfies:

$$
W\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)= \begin{cases}0, & \text { if }|a| \gg 1 \\
|a|^{2}, & \text { if }|a| \ll 1\end{cases}
$$

Hence $W$ is regular $\bmod N$.

Suppose $W^{N} \neq 0$. Then the $\rho\left(G^{1}\right)$-span $X$ of $W^{N}$ is admissible, and every $f$ in $X$ is a function on $G^{1}$ satisfying:

$$
f(n g)=f(g), \quad \text { for all } n \text { in } N \text { and } g \text { in } G^{1}
$$

For each $f$ in $X$, let $\left.f\right|_{A^{1}}$ denote the restriction of $f$ to $A^{1}$. In other words, $\left.f\right|_{A^{1}}$ is the function on $A^{1}$ satisfying:

$$
\left.f\right|_{A^{1}}(a)=f(a), \quad \text { for every } a \text { in } A^{1}
$$

Let $C^{\infty}\left(A^{1}\right)$ denote the space of locally constant functions on $A^{1}$. Note that since each $f$ in $X$ is $K^{1}$-finite, $\left.f\right|_{A^{1}}$ is locally constant on $A^{1}$. So we get a map:

$$
r: X \rightarrow C^{\infty}\left(A^{1}\right) \quad \text { with } r(f)=\left.f\right|_{A^{1}}
$$

Let $X_{N}$ denote the image of $X$ under $r$, and let $\rho_{N}$ denote the natural action of $A^{1}$ on $X_{N}$, given by:

$$
\left.\rho_{N}(a) f\right|_{A^{1}}(b)=f(b a), \quad \text { for all } a, b \text { in } A^{1}
$$

It is clear that $\left(\rho_{N}, X_{N}\right)$ is a smooth $A^{1}$-module.
Claim. $X_{N}$ is an admissible $A^{1}$-module.
Indeed, we may choose (cf. Deligne [2], for example) a neighborhood basis $\left(K_{m}^{1}\right)_{m \geq 1}$ of 1 in $G$ such that
(i) each $K_{m}^{1}$ has an Iwahori factorization: $K_{m}^{1}=N_{m} \cdot A_{m}^{1} \cdot \bar{N}_{m}$, where $N_{m}=N \cap K_{m}^{1}, A_{m}^{1}=A^{1} \cap K_{m}$, and $\bar{N}_{m}=\bar{N} \cap K_{m}^{1}$; and
(ii) $\left(A_{m}^{1}\right)_{m \geq 0}$ is a neighborhood basis of 1 in $A^{1}$ (consisting of compact subgroups). The claim will be proved if we show that $X_{N}^{A_{m}^{1}}$ is finite-dimensional, for every $m \geq 0$.

But for every $m$, the restriction map $r: X \rightarrow X_{N}$ gives rise to maps:

$$
r: X^{K_{m}^{1}} \rightarrow X_{N}^{A_{m}^{1}}, \quad \text { and } \quad r: X^{A_{m}^{1} \overline{N_{m}}} \rightarrow X_{N}^{A_{m}^{1}}
$$

Suppose $f$ in $X$ is fixed under $A_{m}^{1} \overline{N_{m}}$. Then the function $f_{0}$ defined by:

$$
g \mapsto\left(\operatorname{Vol} N_{m}\right)^{-1} \int_{N_{m}}(\rho(n) f)(g) D n
$$

belongs to $X^{K_{m}^{1}}$. Furthermore,

$$
r\left(f_{0}\right)(a)=\left(\operatorname{Vol} N_{m}\right)^{-1} \int_{N_{m}}(\rho(n) f)(a) d n=f(a), \quad \text { for all } a \text { in } A^{1}
$$

Thus the image of $X^{K_{m}^{1}}$ in $X_{N}$ is identical to that of $X^{A_{m}^{1} \overline{N_{m}}}$.
Now let $Z$ be any finite-dimensional subspace of $X_{N}^{A_{m}^{1}}$, and let $Z^{\prime}$ be any finite dimensional subspace of $X$ mapping onto $Z$. Then we can find an open compact subgroup $\bar{O}$ in $\overline{N_{m}}$ such that $Z^{\prime} \subset X^{A_{m}^{l} \bar{O}}$. Choose $a_{0}$ in
$A^{1}$ such that $a_{0}^{-1} \bar{N}_{m} a_{0} \subseteq \bar{O}$. Then if $f$ is in $Z^{\prime}$ and $n$ is in $\bar{N}_{m}$, we have:

$$
\rho(n) \rho\left(a_{0}\right) f=\rho\left(a_{0}\right) \rho\left(a_{0}^{-1} n a_{0}\right) f=\rho\left(a_{0}\right) f .
$$

Thus $\rho\left(a_{0}\right) Z^{\prime} \subseteq X^{A_{m} \bar{N}_{m}}$. And $r\left(\rho\left(a_{0}\right) Z^{\prime}=\rho_{N}\left(a_{0}\right) Z\right.$ is contained in the image of $X^{K_{m}^{1}}$ (by the earlier remark). So $\operatorname{dim} Z=\operatorname{dim} \rho_{N}\left(a_{0}\right) Z \leq$ $\operatorname{dim} X^{K_{m}^{1}}$. This shows that every finite-dimensional subspace of $X_{N}^{A_{m}^{1}}$ has dimension bounded by $\operatorname{dim} X^{K_{m}^{1}}$. Then $X_{N}^{A_{m}^{1}}$ itself is finite-dimensional. This proves the claim.

Now the claim implies in particular that $\left.W^{N}\right|_{A^{1}}$ is $A^{1}$-finite. So there exist a finite number of characters $\mu_{i}$ of $A^{1}$ such that

$$
W^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=\sum_{i=1}^{N} \alpha_{i} \mu_{i}(a), \quad N<\infty, \alpha_{i} \in \mathbf{C} .
$$

But $W^{N}$ : regular implies:

$$
W^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=0, \quad \text { if }|a| \ll 1
$$

This forces each $\alpha_{i}$ to be zero, and $W^{N} \equiv 0$.

## F: archimedean

When $F=\mathbf{C}$, there is nothing to prove due to the absence of discrete series. So assume $F=\mathbf{R}$. The regularity of Whittaker vectors $W$ of square-integrable representations follow from [4], [5]. It remains to show that $W^{N} \equiv 0$.

Let $X$ denote the $\left(\operatorname{Lie} G^{1}, K^{1}\right)$-span of $W^{N}$. Then $(\rho, X)$ is an admissible (Lie $G^{1}, K^{1}$ )-module. As in the $p$-adic case, consider the restriction map $r:\left.f \rightarrow f\right|_{A^{1}}$ from $X$ into the space of smooth functions on $A^{1}$. Let $X_{N}$ denote $\operatorname{Im}(r)$. Then it is stable under the action of $\left(\operatorname{Lie} A^{1}, A_{K}^{1}\right)$. Note that

$$
\rho(n) f(a)=f(a n)=f\left(a n a^{-1} a\right)=f(a),
$$

for every $a$ in $A^{1}$ and $n$ in $N$. So $\rho(n) f-f$ is in $\operatorname{Ker}(r)$ for all $f$ in $X$ and for all $n$ in $N$. On the other hand, if $n=\exp (T)$ with $T$ in $\operatorname{Lie}(N)$, then

$$
\rho(n) f=\rho(\exp T) f=\rho(1+T) f=f+\rho(T) f .
$$

Thus $\rho(n) f-f=\rho(T) f$.
Conversely, for any $T$ in Lie $N, \rho(T) f=r(T-1) f-f$. Hence $\rho($ Lie $N) X=\operatorname{Span}\{\rho(n) f-f\}$. By an earlier remark, we have $\rho(\operatorname{Lie} N) X \subset \operatorname{Ker}(r)$. So $\operatorname{dim} X_{N} \leq \operatorname{dim}(X / \rho(\operatorname{Lie} N) X)$.

Let $\mathfrak{U}($ Lie $N)$ be the univerdal enveloping algebra of (the complexification of) Lie $N$. Then (cf. [1]) the (Lie $G^{1}, K^{1}$ )-module $X$ is finitely generated as an $\mathfrak{U}(\operatorname{Lie} N)$-module. Since

$$
\mathfrak{U}(\operatorname{Lie} N)=\mathbf{C}+\rho(\operatorname{Lie} N) \mathfrak{U}(\operatorname{Lie} N)
$$

we get $\operatorname{dim}(X / \rho(\operatorname{Lie} N) X)<\infty$. Thus $X_{N}$ is finite-dimensional, and hence an admissible ( $\operatorname{Lie} A^{1}, A_{K}^{1}$ )-module.

We can then find a finite set $\left\{\mu_{K} \mid k=1, \ldots, N\right\}$ of characters of Lie $A^{1}$, and scalars $\left\{b_{k}\right\}$ such that

$$
W^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=\sum_{1 \leq k \leq N} b_{k} \mu_{k}\left(\begin{array}{cc}
\log a & 0 \\
0 & -\log a
\end{array}\right)
$$

As the $\mu_{k}$ 's are unitary, there exist real numbers $\left\{t_{k}\right\}$ such that

$$
\mu_{k}\left(\begin{array}{cc}
\log a & 0 \\
0 & -\log a
\end{array}\right)=a^{i t_{k}}
$$

On the other hand, we know, since $W$ is regular, that

$$
W^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) \leq \text { const. }|a|, \quad \text { for }|a| \ll 1
$$

This yields

$$
\left|\sum_{k} a^{i t_{k}}\right|<\text { const. }|a|, \quad \text { for }|a| \ll 1 .
$$

Then for every $j \in\{1,2, \ldots, N\}$

$$
\left|b_{j}+\sum_{k \neq j} b_{k} a^{l\left(t_{k}-t_{j}\right)}\right|<\text { const. }|a| .
$$

As we let $|a| \rightarrow 0$, this forces $b_{J}$ to be zero. This is true for every $j$. Hence

$$
W^{N}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=0, \quad \text { for all } a
$$

So $W^{N}$ is identically zero, and we are done.
REMARK 1.41. $A^{1}$ acts on the characters of $N$ via its conjugation action on $N$. Let $S$ denote the set of equivalence classes of characters of $N$ under this action. Then this set $S$ is finite. Let $\left\{\eta_{1}, \eta_{2}, \ldots, \eta_{r}\right\}$ be a set of representatives for the elements of $S$ in $\hat{N}$. Then it is a standard fact that every discrete series representation $\pi$ of $G^{1}$ admits a Whittaker model relative to some $\eta_{,}, 1 \leq j \leq r$. [4]

We conclude this section with a few remarks on the case of $G=$ GL( $2, F$ ).

## Definition 1.42.

(a) For every $\Phi$ in $S\left(F^{2}\right)$ and $\phi$ in $S\left(F^{*}\right)$, set: $f(g)=f(\Phi, \phi)(g)=f_{\Phi}(g) \cdot \phi(\operatorname{det} g), \quad$ for all $g$ in $G$
(b) $S(N \backslash G)=\left\{f(\Phi, \phi) \mid \Phi \in S\left(F^{2}\right), \phi \in S\left(F^{*}\right)\right\}$
(c) $S_{0}(N \backslash G)=\left\{f(\Phi, \phi) \mid \Phi \in S_{0}\left(F^{2}\right), \phi \in S\left(F^{*}\right)\right\}$.

It may be verified that:
For every $f$ in $S_{0}(N \backslash G)$ and $g$ in $G$,

$$
\begin{equation*}
\theta_{f}(g)=\int_{N} f(w n g) \bar{\eta}(n) d n \tag{1.43}
\end{equation*}
$$

is convergent and defines a $K$-finite function in $L^{2}(N \backslash G, \eta)$.
Clearly we have:

$$
\begin{align*}
\theta_{f(\Phi, \phi)}(g)=\theta_{\Phi}(g) \cdot & (\operatorname{det} g)  \tag{1.44}\\
& \text { for all } \Phi \in S_{0}\left(F^{2}\right), \phi \in S\left(F^{*}\right) \text { and } g \in G
\end{align*}
$$

The following result can be easily established using Th. 1.27:
Proposition 1.45. Suppose $h$ is a $K$-finite function in $L^{2}(N \backslash G, \eta)$ which is orthogonal to $\left\{\theta_{f} \mid f \in S_{0}(N \backslash G)\right\}$. Then for every $t$ in $F^{*}$, the function:

$$
g \mapsto h\left[\left(\begin{array}{cc}
t & 0 \\
0 & 1
\end{array}\right) g\right], \quad g \in G^{1}
$$

is in the discrete part of $L^{2}\left(N \backslash G^{1}, \eta_{t}\right)$, where $\eta_{t}$ denotes the character: $\left(\begin{array}{ll}1 & x \\ 0 & 1\end{array}\right) \mapsto \eta(t x)$ of $N$.
2. The main duality. For any quasi-character $\mu=\left(\mu_{1}, \mu_{2}\right)$ of $A$, and for $h$ in $L^{2}(N \backslash G)$, the formally define the $\mu$-Mellin transform of $h$ by:

$$
\begin{equation*}
L(h, \mu ; g)=\int_{A} h(a g) \delta(a)^{-1 / 2} \mu(a) d^{*} a \quad(\forall g \in G) \tag{2.0}
\end{equation*}
$$

This expression can be seen to be convergent when $\mu$ is a unitary character and $h$ is in $S(N \backslash G)$ or in $\left(\theta\left(S_{0}(N \backslash G)\right)\right)^{N}$.

The result we are after is the following:
Theorem 2.1. For all fin $S_{0}(N \backslash G), \mu$ in $\hat{A}$ and $g$ in $G$,

$$
L(f, \mu ; g)=\frac{1}{2}\left|\varepsilon^{\prime}\left(\mu_{1} \mu_{2}^{-1}\right)\right|^{-2} \cdot L\left(\theta_{f}^{N}, \mu ; g\right)
$$

where ( for any quasi-character $\lambda$ of $F^{*}$ )

$$
\varepsilon^{\prime}(\lambda)=\varepsilon(\lambda, \psi, d x) \cdot \frac{L\left(\alpha \cdot \lambda^{-1}\right)}{L(\lambda)} \quad \text { with } \alpha(x)=|x|
$$

Here $\varepsilon(\lambda, \psi, d x)$ and $L(\lambda)$ are respectively the $\varepsilon$ - and L-factors attached to $\lambda$ (and $\psi$ and $d x$ ) by J. Tate (cf. [13], [14].

Corollary 2.2 ( Preliminary scalar product formula).

$$
\left\langle\theta_{f}, \theta_{h}\right\rangle=\frac{1}{2} \int_{K} d k \int_{\hat{A}} L\left(\theta_{f}^{N}, \mu ; k\right) \overline{L\left(\theta_{h}^{N}, \mu ; k\right)} \cdot\left|\varepsilon^{\prime}\left(\psi_{1} \mu_{2}^{-1}\right)\right|^{-2} d \mu
$$

for all $f, h$ in $S_{0}(N \backslash G)$.
To see that the theorem implies the corollary, note that if $f_{1}, f_{2}$ are in $L^{2}(N \backslash G)$, then by using of the Plancherel formula for $A$, we have:

$$
\left\langle f_{1}, f_{2}\right\rangle=\int_{K \times \hat{A}} L\left(f_{1}, \mu ; k\right) \overline{L\left(f_{2}, \mu ; k\right)} d \mu d k
$$

where $d \mu$ is the dual Haar measure on $\hat{A}$. Hence

$$
\left\langle\theta_{f}, \theta_{h}\right\rangle=\left\langle f, \theta_{h}^{N}\right\rangle=\int_{K} d k \int_{\hat{A}} L(f, \mu ; k) L\left(\theta_{h}^{N}, \mu ; k\right) d \mu
$$

Now we will begin with the proof of the theorem.
Let $\mathbf{B}$ denote the skew-symmetric bilinear form on $F^{2}$ given by $\mathbf{B}(u, v)=a d-b c$, if $u=(a, b)$ and $v=(c, d)$. For $\Phi$ in $S\left(F^{2}\right)$, set:

$$
\begin{equation*}
\hat{\Phi}(v)=\int_{F^{2}} \Phi(u) \psi(\mathbf{B}(u, v)) d u \tag{2.3}
\end{equation*}
$$

The map $\Phi \mapsto \hat{\Phi}$ can be seen to be an automorphism of $S\left(F^{2}\right)$ commuting with the right action of $G^{1}$. It also induces, for each $\sigma$ in $\widehat{K^{1}}$, an automorphism of $S\left(F^{2}, \sigma\right)$ onto itself. Furthermore, if $\Phi$ is in $S\left(F^{2}\right)$, it can be shown for every $L$ in $\mathscr{L}$ with Haar measure $d z$ that $\int_{L} \Phi(z) d z=0$ iff $\int_{L} \hat{\Phi}(z) d z=0$. Hence this symplectic Fourier transform maps $S_{0}\left(F^{2}\right)$ (resp. $S_{0}\left(F^{2}, \sigma\right)$ ) onto itself.

Definition 2.4. For $\Phi$ in $S_{0}\left(F^{2}\right)$ and $y$ in $F^{*}$, set:

$$
\Phi^{*}(y)=\int_{F} \Phi(y, x) d x
$$

Let $F$ be $p$-adic. Then $\Phi^{*}$ clearly vanishes if $|y| \gg 1$. Also, since $\Phi$ is locally constant and since $\int_{F} \Phi(0, x) d x=0$, we have:

$$
\Phi^{*}(y)=0 \quad \text { if }|y| \ll 1
$$

Now assume that $F$ is archimedean. Then $\Phi^{*}(y)$ vanishes rapidly as $|y| \rightarrow \infty$. Moreover, proceeding along the lines of proof of Proposition 1.11, we can show that:

$$
\Phi^{*}(y)<|y|_{F}^{c_{F}}, \quad \text { as }|y| \rightarrow 0
$$

where

$$
c_{\mathbf{R}}=1 \quad \text { and } \quad c_{\mathbf{C}}=\frac{1}{2}
$$

Definition 2.5. For $\Phi$ in $S_{0}\left(F^{2}\right)$, set:

$$
h_{\Phi}(y)= \begin{cases}|y|^{-1} \Phi^{*}(y), & \text { if } y \neq 0 \\ 0, & \text { if } y=0\end{cases}
$$

From the above remarks about the asymptotics of $\Phi$, it can be seen that $h_{\phi}$ is actually an integrable function on $F$.

For any open subset $U$ of $F$, let $S(U)$ denote the space of functions on $U$ which are restrictions to $U$ of functions in $S(F)$. Furthermore, let $U_{\varepsilon}$ denote, for every $\varepsilon>0$, the open subset of $F$ given by: $\{y \in F||y|>\varepsilon\}$.

Definition 2.6.

$$
S_{1}(F)=\left\{h \in L^{1}(F)|h|_{U^{\varepsilon}} \in S\left(U_{\varepsilon}\right), \forall \varepsilon>0\right\}
$$

Then for every $\Phi$ in $S_{0}\left(F^{2}\right), h_{\Phi}$ is in $S_{1}(F)$. Note that when $F$ is $p$-adic, $h_{\Phi}$ is actually in $S(F)$.

For every $h$ in $L^{1}(F)$, let $\hat{h}$ denote the usual Fourier transform relative to $(\psi, d x)$.

Proposition 2.7. Let $f=f(\Phi, \phi)$ be in $S_{0}(N \backslash G)$ (see Definition 1.4.2) with $\Phi \in S_{0}\left(F^{2}\right)$ and $\phi \in S\left(F^{*}\right)$. Then

$$
\theta_{f}^{N}\left(\begin{array}{ll}
a & 0 \\
0 & b
\end{array}\right)=\phi(a b) \cdot|b|^{-1}\left(h_{\Phi}^{\wedge}\left(a^{-1}\right)+\left(h_{\Phi^{\wedge}}\right)^{\wedge}(b)\right)
$$

Given any $f=f(\Phi, \phi)$ in $S_{0}(N \backslash G)$, let $f^{\sim}$ denote the function $f\left(\hat{\Phi}, \Phi^{\vee}\right)$, where $\phi^{\vee}(t)=|t| \cdot \phi\left(t^{-1}\right)$ for $t$ in $F^{*}$.

Corollary 2.8. For every $f$ in $S_{0}(N \backslash G)$ and $g$ in $G$, we have:

$$
\theta_{f}^{N}(g)=\theta_{f^{-}}^{N}\left(g \operatorname{det} \cdot g^{-1}\right)
$$

The corollary easily follows from the proposition when $g$ is a diagonal matrix. We may write any element of $G$ as $n h k$ with $n$ in $N, h$ in $A$ and $k$ in $K^{1}$. So it suffices to show that the truth of the corollary for $g$ implies it for $g k$ as well for $k$ in $K^{1}$. We may take $\Phi$ to be in $S_{0}\left(F^{2}, \sigma\right)$ for some $\sigma$ in $\widehat{K^{1}}$. Then, noting that $\theta_{f}^{N}(g k)$ equals $\sigma\left(k^{-1}\right) \theta_{f}^{N}(g)$, we obtain the desired result.

We also have the following

Corollary 2.9. $f(g)=f^{\sim}\left(g \operatorname{det} g^{-1}\right)$, for all $f$ in $S_{0}(N \backslash G)$ and $g$ in $G$.

Indeed we have

$$
\left\langle f, \boldsymbol{\theta}_{f^{\sim}}^{N}\right\rangle=\left\langle\boldsymbol{\theta}_{f}, \boldsymbol{\theta}_{f^{-}}\right\rangle=\left\langle\boldsymbol{\theta}_{f}^{N}, f^{\sim}\right\rangle
$$

This, together with Corollary 2.8, yields the formula.

Proof of Proposition 2.7. We have to show that

$$
\theta_{f}^{N}\left(\begin{array}{ll}
a & 0 \\
0 & b
\end{array}\right)
$$

equals:

$$
\begin{aligned}
& \phi(a b) \cdot\left[|b|^{-1} \int_{F} \psi\left(a^{-1} y\right)|y|^{-1} d y \int_{F} \Phi(y, x) d x\right. \\
&\left.+|b|^{-1} \int_{F} \psi(b y)|y|^{-1} d y \int_{F} \hat{\Phi}(y, x) d x\right]
\end{aligned}
$$

From the definition of $\theta_{f}$, we get

$$
\theta_{f}^{N}\left(\begin{array}{ll}
a & 0 \\
0 & b
\end{array}\right)=\phi(a b) \int_{F} d x \int_{F} \Phi(a y, b+b x y) \psi(y) d y
$$

Since $\Phi$ is in $S\left(F^{2}\right)$ we can find a function $\Phi_{1}$ which is integrable on $F^{2}$ such that

$$
\Phi(u, v)=\Phi(0, v)+u \Phi_{1}(u, v) \quad \text { for all }(u, v) \text { in } F^{2}
$$

Thus

$$
\theta_{f}^{N}\left(\begin{array}{ll}
a & 0 \\
0 & b
\end{array}\right)=\phi(a b) \cdot(I+J)
$$

where

$$
I=\int d x \int \Phi(0, b+b x y) \psi(y) d y
$$

and

$$
J=\int d x \int a y \Phi_{1}(a y, b+b x y) \psi(y) d y
$$

We write $\int$ for $\int_{F}$. Then

$$
\begin{aligned}
I & =|b|^{-1} \int \frac{d x}{|x|} \int \Phi(0, b+y) \psi\left(b^{-1} x^{-1} y\right) d y, \quad \text { as } y \rightarrow b^{-1} x^{-1} y \\
& =|b|^{-1} \int \frac{d x}{|x|} \int \Phi(0, b+y) \psi(x y) d y, \quad \text { as } x \rightarrow b^{-1} x^{-1} \\
& =|b|^{-1} \int \psi(-b x) \frac{d x}{|x|} \int \phi(0, y) \psi(x, y) d y, \quad \text { as } y \rightarrow y-b \\
& =|b|^{-1} \int \psi(b x) \frac{d x}{|x|} \int \hat{\Phi}(x, y) d y \\
& =|b|^{-1} \int \psi(b y) \frac{d y}{|y|} \int \hat{\Phi}(y, x) d x
\end{aligned}
$$

by simply interchanging the labels $x$ and $y$.
Next we claim that $J$ is absolutely convergent. For this it is enough to show that

$$
J^{*} \stackrel{\text { def }}{=} \int a y \psi(y) d y \int \Phi_{1}(a y, b+b x y) d x
$$

is absolutely convergent. This is indeed so because

$$
\begin{aligned}
J^{*} & =|b|^{-1} \int a y \psi(y) \frac{d y}{|y|} \int \Phi_{1}(a y, b+x) d x, \quad \text { as } x \rightarrow b^{-1} y^{-1} x \\
& =|b|^{-1} \int \psi\left(a^{-1} y\right)(y /|y|) d y \int \Phi_{1}(y, x) d x
\end{aligned}
$$

which is absolutely convergent. So by Fubini's theorem $J=J^{*}$, and we obtain:

$$
\begin{aligned}
J & =|b|^{-1} \int \psi\left(a^{-1} y\right) \frac{d y}{|y|} \int y \Phi_{1}(y, x) d x \\
& =|b|^{-1} \int \psi\left(a^{-1} y\right) \frac{d y}{|y|} \int[\Phi(y, x)-\Phi(0, x)] d x \\
& =|b|^{-1} \int \psi\left(a^{-1} y\right) \int \Phi(y, x) d x, \quad \text { since } \int \Phi(0, x) d x=0
\end{aligned}
$$

For every character $\mu$ of $A$, let $\mu^{w}$ denote the character

$$
\left(\begin{array}{ll}
a & 0 \\
0 & b
\end{array}\right) \mapsto \mu\left(\begin{array}{ll}
b & 0 \\
0 & a
\end{array}\right)
$$

Proposition 2.10. For every fin $S_{0}(N \backslash G)$ and $\mu=\left(\mu_{1}, \mu_{2}\right)$ in $\hat{A}$,

$$
L\left(\theta_{f}^{N}, \mu ; g\right)=\left|\varepsilon^{\prime}\left(\mu_{1} \mu_{2}^{-1}\right)\right|^{2}\left[L(f, \mu ; g)+L\left(f^{\sim}, \mu^{w} ; g \cdot \operatorname{det} g^{-1}\right)\right]
$$

Theorem 2.1 can be easily seen to follow from this proposition and Corollary 2.9.

Replacing $f$ by $\rho(g) f$, we see that we will be done once we show that the proposition holds when $g=e$.

For $h$ in $S_{1}(F)$ and $\nu$ a quasi-character of $F^{*}$, set:

$$
\begin{equation*}
Z(h, \nu)=\int_{F^{*}} h(x) \nu(x) d^{*} x \tag{2.11}
\end{equation*}
$$

This is a local zeta integral as defined by J. Tate ([13]). We have the following functional equation:

$$
\begin{equation*}
Z\left(\hat{h}, \nu^{-1} \cdot \alpha\right)=\varepsilon^{\prime}(\nu) Z(h, \nu) \tag{2.12}
\end{equation*}
$$

where $\alpha(x)=|x|$ for $x$ in $F^{*}$.
A simple calculation, using Proposition 2.7, yields:

$$
\begin{align*}
L\left(\theta_{f}^{N}, \mu ; e\right)= & Z\left(\phi, \mu_{2} \cdot \alpha^{-1 / 2}\right) Z\left(\hat{h}_{\Phi}, \mu_{2} \mu_{1}^{-1}\right)  \tag{2.13}\\
& +Z\left(\phi, \mu_{1} \cdot \alpha^{-1 / 2}\right) Z\left(\hat{g}_{\hat{\Phi}}, \mu_{2} \mu_{1}^{-1}\right)
\end{align*}
$$

By using the functional equation (2.12) and the definition of $h_{\Phi}$ we see that

$$
Z\left(\hat{h}_{\Phi}, \mu_{2} \mu_{1}^{-1}\right)=\varepsilon^{\prime}\left(\mu_{2}^{-1} \mu_{1} \alpha\right) \int_{F^{*} \times F} \Phi(a, x) \mu_{2}^{-1} \mu^{1}(a) d^{*} a d x
$$

which equals

$$
\left[\frac{\varepsilon^{\prime}\left(\mu_{1} \mu_{2}^{-1} \alpha\right)}{\varepsilon^{\prime}\left(\mu_{1} \mu_{2}^{-1}\right)}\right] \int_{F^{*}} \hat{\Phi}(0, b) \mu_{1} \mu_{2}^{-1}(b)|b| d^{*} b
$$

From this it can be deduced that

$$
\begin{equation*}
L(f, \mu ; e)=\left[\frac{\varepsilon^{\prime}\left(\mu \mu_{2}^{-1} \alpha\right)}{\varepsilon^{\prime}\left(\mu_{1} \mu_{2}^{-1}\right)}\right] \cdot Z\left(\phi, \mu_{1} \alpha^{-1 / 2}\right) \cdot Z\left(\hat{h}_{\Phi}, \mu_{2} \mu_{1}^{-1}\right) \tag{2.14}
\end{equation*}
$$

Similarly, we get

$$
\begin{equation*}
L\left(f^{\sim}, \mu^{w} ; e\right)=\left[\frac{\varepsilon^{\prime}\left(\mu_{1} \mu_{2}^{-1} \alpha\right)}{\varepsilon^{\prime}\left(\mu_{1} \mu_{2}^{-1}\right)}\right] \cdot Z\left(\phi, \mu_{2} \alpha^{-1 / 2}\right) \cdot Z\left(\hat{h}_{\Phi}, \mu_{2} \mu_{1}^{-1}\right) \tag{2.15}
\end{equation*}
$$

The proposition, and hence the theorem, follow from (2.13), (2.14), and (2.15), once we note (cf. [13]) that when $\nu$ is a character (of module one) of $F^{*}$,

$$
\left[\varepsilon^{\prime}(\nu) / \varepsilon^{\prime}(\nu \alpha)\right]=\left|\varepsilon^{\prime}(\nu)\right|^{2}
$$

3. The scalar product formula. For any $\Phi$ in $S\left(F^{2}\right)$, not necessarily in $S_{0}\left(F^{2}\right)$, and quasicharacter $\mu=\left(\mu_{1}, \mu_{2}\right)$ of $A$, set (as in [7]):

$$
\begin{align*}
& b_{\Phi}(\mu ; g)  \tag{3.0}\\
& \quad=\mu_{1}(\operatorname{det} g)|\operatorname{det} g| \int_{F^{*}} \Phi[(0, t) g] \mu_{1} \mu_{2}^{-1}(t)|t| d^{*} t, \quad \forall g \in G
\end{align*}
$$

Then $b_{\Phi}(\mu)$ belongs to the space of the smoothly induced representation $\operatorname{Ind}(G, B ; \mu)$, and every function in this space is of this form. The associated Whittaker vector is:

$$
\begin{equation*}
W_{\Phi}(\mu ; g)=\int_{N} b_{\Phi}(\mu ; w n g) \overline{\eta(n)} d n \quad(g \in G) \tag{3.1}
\end{equation*}
$$

For any $\xi$ in $L^{2}(N \backslash G, \eta)$, formally set:

$$
\begin{equation*}
\Psi\left(\xi, W_{\Phi}(\mu)\right)=\int_{M \backslash G} \xi(g) W_{\Phi}(\mu ; \delta g) d g \tag{3.2}
\end{equation*}
$$

where $\delta=\left(\begin{array}{rl}-1 & 0 \\ 0 & 1\end{array}\right)$.
It may be seen that, when $\xi=\theta_{f}$ for some $f$ in $S(N \backslash G)$ and when $W_{\Phi}(\mu)$ belongs to a unitary principal series representation, this expression makes sense. Furthermore it equals, in the sense of analytic continuation:

$$
\begin{equation*}
\mu\left(w \delta w^{-1}\right) \int_{K} b_{\Phi}(\mu ; k) L\left(\xi^{N}, \mu ; k\right) d k \tag{3.3}
\end{equation*}
$$

Definition 3.4. For every $\sigma$ in $\hat{K}$, let $I_{\sigma}$ denote the set of pairs $(i, j)$ of integers such that $1 \leq i, j \leq \operatorname{dim} \sigma$, and write $\Lambda=\left\{\left(\sigma, I_{\sigma}\right) \mid \sigma \in \hat{K}\right\}$.

Theorem 3.5. There exists a family $\left\{\Phi_{l} \in S\left(F^{2}\right) \mid l \in \Lambda\right\}$ such that, for all $f, h$ in $S_{0}(N \backslash G)$,

$$
\left\langle\theta_{f}, \theta_{h}\right\rangle=\frac{1}{2} \sum_{l \in \Lambda} \cdot \int_{\hat{A}} \Psi\left(\theta_{f}, W_{\Phi_{l}}(\mu)\right) \cdot \overline{\Psi\left(\theta_{h}, W_{\Phi_{l}}(\mu)\right)} \cdot\left|\varepsilon^{\prime}\left(\mu_{1} \mu_{2}^{-1}\right)\right|^{-2} d \mu
$$

Proof. Let $\left(\sigma, H_{\sigma}\right)$ be a unitary irreducible of $K^{1}$, and let $\mu$ be a character of $A$. Then, for $\beta$ in $L^{2}(N \backslash G)$, we can formally define an
operator:

$$
\begin{equation*}
L(\beta, \mu ; \sigma)=\iint_{A K} \sigma(k) \beta(a k) \mu(a) \delta(a)^{-1 / 2} d^{*} a d k \tag{3.6}
\end{equation*}
$$

This expression converges whenever $\beta$ is in $S(N \backslash G)$ or in $\left(\theta\left(S_{0}(N \backslash G)\right)\right)^{N}$.

Let $L(\beta, \mu ; \sigma)^{*}$ be the ajoint of $L(\beta, \mu ; \sigma)$. Then, for all $f, h$ in $S_{0}(N \backslash G)$, the preliminary scalar product formula (Corollary 2.2) yields:

$$
\begin{equation*}
\left\langle\theta_{f}, \theta_{h}\right\rangle=\frac{1}{2} \sum_{\hat{K}} \int_{\hat{A}} \operatorname{sp}\left[L\left(\theta_{h}^{N}, \mu ; \sigma\right)^{*} L\left(\theta_{f}^{N}, \mu ; \sigma\right)\right]\left|\varepsilon^{\prime}\left(\mu_{1}^{-1} \mu_{2}\right)\right|^{-2} d \mu \tag{3.7}
\end{equation*}
$$

where sp denotes tr $/ \operatorname{deg} \sigma$.
For each character $\mu$ of $A$, let $\mu_{0}$ denote the restriction to $A_{K}$. Recall that the eigenspace $H_{\sigma}\left(\mu_{0}\right)$ is at most one-dimensional (cf. [3]). Let $P\left(\sigma, \mu_{0}\right)$ be the corresponding projector. Then

$$
\begin{equation*}
P\left(\sigma, \mu_{0}\right) L(\beta, \mu ; \sigma)=L(\beta, \mu ; \sigma) \tag{3.8}
\end{equation*}
$$

When $F$ is $p$-adic, let $\chi$ denote the characteristic function of the unit group $U_{F}$ divided by the volume of $U_{F}$. For $t$ in $F^{*}$, set:

$$
\chi_{\mu}(t)= \begin{cases}\mu(t) \cdot \chi(t), & \text { if } F: p \text {-adic }  \tag{3.9}\\ \mu(t) \cdot \exp \left(-\pi t^{2}\right), & \text { if } F=\mathbf{R} \\ \mu(t) \cdot \exp (-2 \pi t \bar{t}), & \text { if } F=\mathbf{C}\end{cases}
$$

By construction $\int_{F^{*}} \chi_{\mu}(t) \mu^{-1}(t)|t| d^{*} t=1$. Set

$$
\begin{equation*}
\Phi_{\sigma}[(0, t) k]=\chi_{\mu}(t) \sigma\left(k^{-1}\right) P\left(\sigma, \mu_{0}\right) \tag{3.10}
\end{equation*}
$$

This defines an operator-valued function $\Phi_{\sigma}$ in $S\left(F^{2}\right)$. Furthermore,

$$
b_{\Phi_{o}}(\bar{\mu} ; \delta)=\int_{F^{*}} \Phi_{\sigma}(0, t) \mu^{-1}(t) d^{*} t=P\left(\sigma, \mu_{0}\right)
$$

Thus we get, using (3.8)

$$
\begin{equation*}
L(\beta, \mu ; \sigma)=\int_{N \backslash G} \beta(g) b_{\Phi_{o}}\left(\bar{\mu}^{w} ; g\right) d g \tag{3.11}
\end{equation*}
$$

From (3.7) and (3.11) we get:

$$
\begin{align*}
&\left\langle\phi_{f}, \theta_{h}\right\rangle=\frac{1}{2} \sum_{\sigma \in \hat{K}^{\prime}} \int_{\hat{A}} d \mu \iint_{(N G)^{2}}\left(\theta_{f}^{N}(g) \overline{\theta_{h}^{N}\left(g^{\prime}\right)} \cdot\left|\varepsilon^{\prime}\left(\mu_{1}^{-1} \mu_{2}\right)\right|^{-2}\right)  \tag{3.12}\\
& \cdot \operatorname{sp}\left[b_{\Phi_{\sigma}}\left(\bar{\mu}^{w} ; g\right)^{*} \cdot b_{\Phi_{\sigma}}\left(\bar{\mu}^{w} ; g^{\prime}\right)\right] d g d g^{\prime}
\end{align*}
$$

Finally, choosing a basis $\left\{e_{j}\right\}$ of each $H_{\sigma}$, set

$$
\begin{equation*}
\Phi_{i, j}=\left(\operatorname{dim} H_{\sigma}\right)^{-1} \cdot\left\langle\Phi e_{i}, \Phi e_{j}\right\rangle \tag{3.13}
\end{equation*}
$$

Then

$$
\begin{align*}
& \operatorname{sp}\left[b_{\Phi_{\sigma}}\left(\mu^{w} ; g\right)^{*} b_{\Phi_{\sigma}}\left(\mu^{w} ; g\right)^{*} b_{\Phi_{o}}\left(\mu^{w} ; g^{\prime}\right)\right]  \tag{3.14}\\
&=\sum_{1 \leq i, j \leq \operatorname{dim} H_{\sigma}} b_{\Phi_{i j}}\left(\mu^{w} s ; g^{\prime}\right) b_{\Phi_{t, j}}\left(\mu^{w} ; g\right)
\end{align*}
$$

Now the theorem follows from (3.2), (3.7), (3.12) and (3.14).
Note added in proof. The main results of this paper admit an extension to $\mathrm{GL}_{n}$ by an analogous but not so explicit a method ([12]). By employing different (and much more powerful) principles involving the weak inequality and Maass-Selberg relations, Professor Harish Chandra obtained a complete spectral decomposition in 1982 valid for all real reductive groups. A sketch of proof is in his unpublished paper, "On the theory of the Whittaker integral."

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