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# THE SPIN L-FUNCTION OF QUASI-SPLIT $D_{4}$ 

WEE TECK GAN AND JOSEPH HUNDLEY

## 1. Introduction

The purpose of this paper is to construct a Rankin-Selberg integral for the Spin L-function $L\left(\underline{s}, \pi, V_{\mathrm{Spin}, E}\right)$ of a generic cuspidal representation $\pi$ of a quasi-split adjoint group $G_{E}$ of type $D_{4}$. Let us formulate our results more precisely.
1.1. Étale cubic algebras. Let $F$ be a number field with adele ring $\mathbb{A}$ and absolute Galois $\operatorname{group} \operatorname{Gal}(\bar{F} / F)$. An étale cubic algebra is an $F$-algebra $E$ such that $E \otimes_{F} \bar{F} \cong \bar{F}^{3}$. More concretely, an étale cubic $F$-algebra is of the form:

$$
E=\left\{\begin{array}{l}
F \times F \times F ; \\
F \times K, \text { where } K \text { is a quadratic field extension of } F ; \\
\text { a cubic field. }
\end{array}\right.
$$

Since the split algebra $F \times F \times F$ has automorphism group $S_{3}$ (the symmetric group on 3 letters), the isomorphism classes of étale cubic algebras $E$ over $F$ are naturally classified by the set of conjugacy classes of homomorphisms

$$
\rho_{E}: \operatorname{Gal}(\bar{F} / F) \longrightarrow S_{3} .
$$

By composing the homomorphism $\rho_{E}$ with the sign character of $S_{3}$, we obtain a quadratic character (possibly trivial) of $\operatorname{Gal}(\bar{F} / F)$ which corresponds to an étale quadratic algebra $K_{E}$. We call $K_{E}$ the discriminant algebra of $E$. To be concrete,

$$
K_{E}=\left\{\begin{array}{l}
F \times F, \text { if } E=F^{3} \text { or a cyclic cubic field; } \\
K, \text { if } E=F \times K \\
\text { the unique quadratic subfield in the Galois closure of } E \text { otherwise. }
\end{array}\right.
$$

We shall let $\chi_{K_{E}}$ denote the quadratic idele class character associated to $K_{E}$.
1.2. Twisted form of $S_{3}$. Fix an étale cubic $F$-algebra $E$. Then, via the associated homomorphism $\rho_{E}, \operatorname{Gal}(\bar{F} / F)$ acts on $S_{3}$ (by inner automorphisms) and thus defines a twisted form $\Gamma_{E}$ of the finite constant group scheme $S_{3}$. For any commutative $F$-algebra $A$, we have

$$
\Gamma_{E}(A)=\operatorname{Aut}_{A}\left(E \otimes_{F} A\right)
$$

1.3. Quasi-split groups of type $D_{4}$. Because $S_{3}$ is the outer automorphism group of $\mathrm{PGSO}_{8}$ (the split adjoint group of type $D_{4}$ ), associated to $E$ is a quasi-split adjoint group $G_{E}$ of type $D_{4}$. The outer automorphism group of $G_{E}$ is precisely the finite group scheme $\Gamma_{E}$.

The Langlands dual group of $G_{E}$ is the simply-connected complex Lie group

$$
G_{E}^{\vee}=\operatorname{Spin}_{8}(\mathbb{C})
$$

This has 3 irreducible 8-dimensional representations $V_{\mu_{i}}$, whose highest weights $\mu_{i}(i=1,2$ or 3) are the three fundamental weights associated to the three satellite vertices in the Dynkin diagram of type $D_{4}$. These 3 representations are permuted by the outer automorphism group $S_{3}$, and the sum

$$
V_{S p i n}=V_{\mu_{1}} \oplus V_{\mu_{2}} \oplus V_{\mu_{3}}
$$

extends to a faithful irreducible representation of $\operatorname{Spin}_{8}(\mathbb{C}) \rtimes S_{3}$. In fact, there are two such extensions which differ from each other by twisting with the sign character of $S_{3}$. So we shall need to specify the extension precisely later (in (2.8)).

The L-group ${ }^{L} G_{E}$ is a certain semidirect product of $\operatorname{Spin}_{8}(\mathbb{C})$ with $\operatorname{Gal}(\bar{F} / F)$. More precisely, the action of $\operatorname{Gal}(\bar{F} / F)$ on $\operatorname{Spin}_{8}(\mathbb{C})$ is via the homomorphism $\rho_{E}$. Thus there is a natural map

$$
{ }^{L} G_{E} \longrightarrow \operatorname{Spin}_{8}(\mathbb{C}) \rtimes S_{3}
$$

whose restriction to $\operatorname{Gal}(\bar{F} / F)$ is $\rho_{E}$. Via this map, we may view $V_{S p i n}$ as a representation of ${ }^{L} G_{E}$. We denote this representation by $V_{\text {Spin }, E}$ and call this the Spin representation of ${ }^{L} G_{E}$. Note that $V_{\text {Spin, } E}$ is reducible unless $E$ is a field.
1.4. The Spin L-function. Now let $\pi=\otimes_{v} \pi_{v}$ be a cuspidal automorphic representation of $G_{E}(\mathbb{A})$. For almost all $v$, the representation $\pi_{v}$ is unramifed and gives rise to a semisimple conjugacy class $t_{\pi_{v}} \in{ }^{L} G_{E}$ (its Satake parameter). We may thus define the partial Spin L-function associated to $\pi$ :

$$
L^{S}\left(s, \pi, V_{\mathrm{Spin}, E}\right)=\prod_{v \notin S} \frac{1}{\operatorname{det}\left(1-q_{v}^{-s} t_{\pi_{v}} \mid V_{\mathrm{Spin}, E}\right)}
$$

where $S$ is a finite set of places, including the archimedean ones, such that $\pi_{v}$ is unramified for all $v \notin S$. This Euler product converges when $\operatorname{Re}(s)$ is sufficiently large.

In fact, since the representation $V_{\text {Spin, } E}$ is reducible when $E$ is not a field, it is natural to define a refinement of the above $L$-function by introducing a different variable for each irreducible constituent of $V_{\text {Spin }, E}$. More precisely, we have:

- If $E=F \times F \times F$, we set $\underline{s}=\left(s_{1}, s_{2}, s_{3}\right)$ and

$$
L^{S}\left(\underline{s}, \pi, V_{\text {Spin }, E}\right)=\prod_{v \notin S} \frac{1}{\operatorname{det}\left(1-q_{v}^{-s_{1}} t_{\pi_{v}} \mid V_{\mu_{1}}\right) \cdot \operatorname{det}\left(1-q_{v}^{-s_{2}} t_{\pi_{v}} \mid V_{\mu_{2}}\right) \cdot \operatorname{det}\left(1-q_{v}^{-s_{3}} t_{\pi_{v}} \mid V_{\mu_{3}}\right)} .
$$

- If $E=F \times K$, we may assume without loss of generality that the two fundamental weights $\mu_{2}$ and $\mu_{3}$ are permuted by the Galois action, so that $V_{\mu_{1}}$ and $V_{\mu_{2}} \oplus V_{\mu_{3}}$ are the two irreducible constituents of $V_{\text {Spin, } E}$. Setting $\underline{s}=\left(s_{1}, s_{23}, s_{23}\right)$, we then define

$$
L^{S}\left(\underline{s}, \pi, V_{\mathrm{Spin}, E}\right)=\prod_{v \notin S} \frac{1}{\operatorname{det}\left(1-q_{v}^{-s_{1}} t_{\pi_{v}} \mid V_{\mu_{1}}\right) \cdot \operatorname{det}\left(1-q_{v}^{-s_{23}} \pi_{\pi_{v}} \mid V_{\mu_{2}} \oplus V_{\mu_{3}}\right)} .
$$

- if $E$ is a field, then $\underline{s}=(s, s, s)$ and $L^{S}\left(\underline{s}, \pi, V_{\mathrm{Spin}, E}\right)$ is as originally defined.
1.5. The refined zeta function of $E$. Likewise, since $E$ is the product of fields $E_{i}$, the zeta function $\zeta_{E}(s)$ decomposes as the product of zeta functions of these fields. It is natural to define a refinement $\zeta_{E}(\underline{s})$, with $\underline{s}$ defined as above in the three different cases. Thus, for example,

$$
\zeta_{F \times K}(\underline{s})=\underset{2}{\zeta_{F}\left(s_{1}\right) \cdot \zeta_{K}\left(s_{23}\right) .}
$$

In addition, we define another refinement of the zeta function of $E$ :

$$
\hat{\zeta}_{E}(\underline{s})=\left\{\begin{array}{l}
\zeta\left(s_{1}+s_{2}-s_{3}\right) \cdot \zeta\left(s_{2}+s_{3}-s_{1}\right) \cdot \zeta\left(s_{3}+s_{1}-s_{2}\right), \text { if } E=F^{3} \\
\zeta\left(2 s_{23}-s_{1}\right) \cdot \zeta_{K}\left(s_{1}\right), \text { if } E=F \times K \\
\zeta_{E}(s), \text { if } E \text { is a field. }
\end{array}\right.
$$

This may seem somewhat artificial at first, but note that if one sets all the variables $s_{i}$ equal to $s$, one recovers $\zeta_{E}(s)$. In fact, this refined function arises naturally in the normalizing factor of the Eisenstein series in our Rankin-Selberg integral. In any case, one can simply regard it as a shorthand for the rather convoluted function it represents.
1.6. The Rankin-Selberg integral. The goal of this paper is to construct a Rankin-Selberg integral for $L\left(\underline{s}, \pi, V_{\text {Spin }, E}\right)$ when $\pi$ is globally generic, i.e. possesses a non-zero WhittakerFourier coefficient. The Rankin-Selberg integral we consider is of Shimura type. Thus it involves the integral of a cusp form, an Eisenstein series and a "theta" function:

$$
Z_{E}(\varphi, \Phi, f, \underline{s})=\int_{G_{E}(F) \backslash G_{E}(\mathbb{A})} \varphi(g) \cdot \overline{\theta(f)(g)} \cdot E(\Phi, \underline{s}, g) d g
$$

To explain the various notations,

- $\varphi$ is a cusp form in $\pi$;
- $\theta(f)$ is a vector in the minimal representation $\Pi_{E}$ of $G_{E}(\mathbb{A})$, which is the analog of the Weil representation of a metaplectic group;
- $E(\Phi, \underline{s}, g)$ is an Eisenstein series on $G_{E}(\mathbb{A})$ associated to a standard section $\Phi_{\underline{s}}$ of a certain principal series $I_{Q_{E}}(\underline{s})$ of $G_{E}$. This principal series $I_{Q_{E}}(\underline{s})$ is induced from a character on the (not-necessarily maximal) parabolic subgroup $Q_{E}$ whose Levi factor is of semisimple type $A_{1}$ corresponding to the middle vertex of the Dynkin diagram of type $D_{4}$.

Our main theorem is then:
1.7. MAIN THEOREM: The global zeta integral admits an Euler product. Moreover, we have:

$$
Z_{E}(\varphi, \Phi, f, \underline{s})=\frac{L^{S}\left(\underline{s}, \pi, V_{\mathrm{Spin}, E}\right)}{\hat{\zeta}_{E}^{S}(\underline{s}+\underline{1}) \cdot \zeta_{E}^{S}(2 \underline{s}) \cdot \zeta^{S}(|\underline{s}|)} \cdot\left(\prod_{v \in S} Z_{E, v}\left(\varphi_{v}, \Phi_{v}, f_{v}, \underline{s}\right)\right)
$$

where $\underline{1}=(1, . ., 1)$ and $|\underline{s}|$ denotes the sum of the components of $\underline{s}$ (so for example, $|\underline{s}|=3 s$ if $E$ is a field). Moreover, each local factor for $v \in S$ admits meromorphic continuation as a function of $\underline{s}$. Thus, the Spin L-function $L^{S}\left(\underline{s}, \pi, V_{\text {Spin, }}\right)$ can be meromorphically continued.

Our investigation is inspired by the recent paper [GH] of the second author with D. Ginzburg. There, they considered the case when the group $G_{E}$ is split and constructed a multi-variable Rankin-Selberg integral which is inherently asymmetric: it gives the $L$ functions associated to two of the degree 8 representations (say $V_{\mu_{1}}$ and $V_{\mu_{2}}$ ), with one of them obtained two times (say $V_{\mu_{1}}$ ). The Rankin-Selberg integral we consider here is motivated by the $S_{3}$-symmetry of $D_{4}$, and has the advantage that it extends in a self evident fashion to all the quasi-split forms. More importantly, when $E$ is a field, the Spin $L$-function is one which cannot be analyzed by the Langlands-Shahidi method.

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## 2. Preliminaries

We begin by establishing some notations and introducing some background results.
2.1. $D_{4}$ root system. We fix a Borel subgroup $B_{E}$ of $G_{E}$ and a Levi subgroup $T_{E}$ which is a maximal torus, both defined over $F$. Let $V_{E}$ denote the unipotent radical of $B_{E}$. Over the algebraic closure $\bar{F}$, the choice of $\left(T_{E}, B_{E}\right)$ gives a set $\Delta$ of simple absolute roots for the set $\Psi$ of absolute roots of $T_{E}$ in $G_{E}$. Using the standard realization of the root system of type $D_{4}$, we label these simple roots by

$$
\left\{\begin{array}{l}
\alpha_{1}=\epsilon_{1}-\epsilon_{2} \\
\alpha_{0}=\epsilon_{2}-\epsilon_{3} \\
\alpha_{2}=\epsilon_{3}-\epsilon_{4} \\
\alpha_{3}=\epsilon_{3}+\epsilon_{4},
\end{array}\right.
$$

where $\left\{\epsilon_{i}\right\}$ is the standard basis of $\mathbb{R}^{4}$. In particular, $\alpha_{0}$ is the branch (or middle) vertex in the Dynkin diagram of type $D_{4}$. We let

$$
\beta_{0}=\alpha_{1}+\alpha_{2}+\alpha_{3}+2 \alpha_{0}
$$

denote the highest root in $\Psi$. Here is the Dynkin diagram:


For each absolute root $\gamma$, we let $U_{\gamma}$ be the associated root subgroup; it may not be defined over $F$. Indeed, the absolute Galois group $\operatorname{Gal}(\bar{F} / F)$ acts naturally on $\Psi$ and $U_{\gamma}$ is defined over $F$ iff $\gamma$ is fixed by the Galois action. We also let $w_{\gamma}$ denote the element of the Weyl group corresponding to the reflection in $\gamma$.

If $E$ is a cubic field, then $\operatorname{Gal}(\bar{F} / F)$ permutes the roots $\alpha_{1}, \alpha_{2}$ and $\alpha_{3}$ transitively. If $E=F \times K$ with $K$ a quadratic field, then without loss of generality, we assume that $\alpha_{1}$ is fixed, whereas $\alpha_{2}$ and $\alpha_{3}$ are exchanged by the Galois action. If $E$ is the split algebra, the Galois action on $\Psi$ is trivial.

Let $\omega_{i}$ be the fundamental weight so that

$$
\left\langle\omega_{i}, \alpha_{j}^{\vee}\right\rangle=\delta_{i j},
$$

and let $\mu_{i}$ be the fundamental coweight so that

$$
\left\langle\alpha_{j}, \mu_{i}\right\rangle=\delta_{i j} .
$$

In terms of the basis $\left\{\epsilon_{i}^{*}\right\}$ dual to $\left\{\epsilon_{i}\right\}$, the fundamental coweights are given by:

$$
\left\{\begin{array}{l}
\mu_{1}=\epsilon_{1}^{*}, \\
\mu_{0}=\epsilon_{1}^{*}+\epsilon_{2}^{*} \\
\mu_{2}=\left(\epsilon_{1}^{*}+\epsilon_{2}^{*}+\epsilon_{3}^{*}-\epsilon_{4}^{*}\right) / 2 \\
\mu_{3}=\left(\epsilon_{1}^{*}+\epsilon_{2}^{*}+\epsilon_{3}^{*}+\epsilon_{4}^{*}\right) / 2 .
\end{array}\right.
$$

Because $G_{E}$ is adjoint, the set $\left\{\mu_{i}\right\}$ is a basis for the cocharacter group

$$
X_{*}\left(T_{E}\right)=\operatorname{Hom}_{\bar{F}}\left(\mathbb{G}_{m}, T_{E}\right) .
$$

2.2. $G_{2}$ root system. We have already fixed the pair $\left(T_{E}, B_{E}\right)$. If we further fix a ChevalleySteinberg system of épinglage relative to this pair, then we have a compatible system of isomorphisms $U_{\gamma} \cong \mathbb{G}_{a}$ defined over $\bar{F}$ which are permuted by $\operatorname{Gal}(\bar{F} / F)$. This gives a splitting of the outer automorphism group

$$
\Gamma_{E} \hookrightarrow \operatorname{Aut}\left(G_{E}\right) .
$$

The subgroup scheme of $G_{E}$ fixed pointwise by $\Gamma_{E}$ is independent of the choice of the épinglage relative to ( $T_{E}, B_{E}$ ), and is isomorphic to the split exceptional group of type $G_{2}$.

Observe that $B=G_{2} \cap B_{E}$ is a Borel subgroup of $G_{2}$ and $T=T_{E} \cap G_{2}$ is a maximal split torus of $G_{2}$. The torus $T$ is such that

$$
X_{*}(T)=\left\langle\mu_{0}, \mu_{1}+\mu_{2}+\mu_{3}\right\rangle .
$$

Via the adjoint action of $T$ on $G_{E}$, we obtain the root system $\Psi_{G_{2}}$ of $G_{2}$, so that

$$
\Psi_{G_{2}}=\left.\Psi\right|_{T}
$$

Here is a diagram of the root system of type $G_{2}$.


We denote the short simple root of this $G_{2}$ root system by $\alpha$ and the long simple root by $\beta$. Then

$$
\beta=\left.\alpha_{0}\right|_{T} \quad \text { and } \quad \alpha=\left.\alpha_{1}\right|_{T}=\left.\alpha_{2}\right|_{T}=\left.\alpha_{3}\right|_{T} .
$$

Thus, the short root spaces have dimension 3, whereas the long root spaces have dimension 1. For each root $\gamma \in \Psi_{G_{2}}$, the associated root subgroup $U_{\gamma}$ is defined over $F$ and the Chevalley-Steinberg system of épinglage gives isomorphisms:

$$
U_{\gamma} \cong\left\{\begin{array}{l}
\operatorname{Res}_{E / F} \mathbb{G}_{a}, \text { if } \gamma \text { is short; } \\
\mathbb{G}_{a}, \text { if } \gamma \text { is long. }
\end{array}\right.
$$

When $E$ is a cubic field, $T$ is in fact the maximal $F$-split torus of $G_{E}$ and $\Psi_{G_{2}}$ is the relative root system of $G_{E}$.

For each $\gamma \in \Psi_{G_{2}}$, we shall also let $N_{\gamma}$ denote the root subgroup of $G_{2}$ corresponding to $\gamma$. In particular,

$$
N_{\gamma}=U_{\gamma} \cap G_{2}
$$

Because the highest root $\beta_{0}$ of the $D_{4}$-root system restricts to that of the $G_{2}$-root system, we shall let $\beta_{0}$ denote the highest root of the $G_{2}$-root system also. This should not cause confusion.
2.3. Two parabolic subgroups. The $G_{2}$ root system gives rise to 2 parabolic subgroups of $G_{E}$. One of these is a maximal parabolic $P_{E}=M_{E} N_{E}$ known as the Heisenberg parabolic. Its unipotent radical $N_{E}$ is a Heisenberg group and its Levi subgroup $M_{E}$ is spanned by the 3 satellite vertices in the Dynkin diagram. The other parabolic $Q_{E}=L_{E} U_{E}$ is a not-necessarily-maximal parabolic; its Levi subgroup $L_{E}$ is spanned by the branch vertex $\alpha_{0}$ and its unipotent radical $U_{E}$ is a 3 -step unipotent group. We shall need to examine the structure of these 2 parabolics more carefully.
2.4. The Heisenberg parabolic $P_{E}$. Let us begin with the Heisenberg parabolic $P_{E}=$ $M_{E} N_{E}$. Its unipotent radical is a Heisenberg group with center $Z=U_{\beta_{0}}$. Moreover,

$$
N_{E} / Z=U_{\beta} \times U_{\beta+\alpha} \times U_{\beta+2 \alpha} \times U_{\beta+3 \alpha} \cong F \times E \times E \times F .
$$

Note that $P_{E} \cap G_{2}$ is the Heisenberg maximal parabolic $P=M N$ of $G_{2}$, with

$$
M=G_{2} \cap M_{E} \cong G L_{2} \quad \text { and } \quad N=G_{2} \cap N_{E} .
$$

Let $\Omega_{E}(F)$ denote the minimal non-zero $M_{E}(F)$-orbit on

$$
\bar{N}_{E}(F) / \bar{Z}(F)=U_{-\beta} \times U_{-\beta-\alpha} \times U_{-\beta-2 \alpha} \times U_{-\beta-3 \alpha} \cong F \times E \times E \times F .
$$

It is the orbit of a highest weight vector and its Zariski closure is a cone. A non-zero element $(a, x, y, d)$ lies in $\Omega_{E}$ iff

$$
\left\{\begin{array}{l}
x y=a d \\
x^{\#}=a y \\
y^{\#}=d x .
\end{array}\right.
$$

Here

$$
\#: E \longrightarrow E
$$

is the canonical quadratic map with the property that $x \cdot x^{\#}=\mathbb{N}(x)$, where $\mathbb{N}(x)$ denotes the norm of $x$. So, for example, the element $a \cdot(1, x, y, d)$, with $a \in F^{\times}$, lies in $\Omega_{E}(F)$ iff

$$
y=x^{\#} \quad \text { and } \quad d=\mathbb{N}(x) .
$$

Observe that there is a natural map $\bar{N}_{E} / \bar{Z} \longrightarrow \bar{N} / \bar{Z}$ given by

$$
(a, x, y, d) \mapsto(a, \operatorname{Tr}(x), \operatorname{Tr}(y), d)
$$

where $\operatorname{Tr}(x)$ denotes the trace of $x$. Given any element $\chi \in \bar{N} / \bar{Z}$, we let $\Omega_{E, \chi}$ denote the fiber of this map over $\chi$. For example, if $E$ is a field and $\chi=(1,0,0,0)$, then $\Omega_{E, \chi}=\{(1,0,0,0)\}$ since the only $x \in E$ with $\mathbb{N}(x)=0$ is $x=0$.
2.5. The 3-step parabolic $Q_{E}$. Now we come to the parabolic $Q_{E}$. The unipotent radical $U_{E}$ has a filtration

$$
\{1\} \subset U_{E}^{(1)} \subset U_{E}^{(2)} \subset U_{E}
$$

such that

$$
U_{E}^{(1)}=U_{\beta_{0}} \times U_{\beta_{0}-\beta}
$$

is the center of $U_{E}$. Further,

$$
U_{E}^{(2)}=\left[U_{E}, U_{E}\right]=U_{\beta_{0}} \times U_{\beta_{0}-\beta} \times U_{2 \alpha+\beta}
$$

is the commutator subgroup of $U_{E}$ and is abelian. In particular, $U_{E}$ is a 3 -step unipotent group; hence we call $Q_{E}$ the 3 -step parabolic. Note that $Q=G_{2} \cap Q_{E}=L \cdot U$ is the non-Heisenberg maximal parabolic of $G_{2}$, with

$$
L=G_{2} \cap L_{E} \cong G L_{2} \text { and } \quad U=G_{2} \cap U_{E}
$$

It will be necessary to have a more explicit description of $L_{E}$. By examining the root datum of $L_{E}$, one can show that

$$
L_{E} \cong\left(G L_{2} \times \operatorname{Res}_{E / F} \mathbb{G}_{m}\right) / \Delta \mathbb{G}_{m}
$$

In terms of this isomorphism, the simple absolute roots of $G_{E}$ are given on a torus element

$$
t=\left(\begin{array}{ll}
a & 0 \\
0 & b
\end{array}\right) \times\left(c_{1}, c_{2}, c_{3}\right) \in T_{E}(\bar{F}) \subset L_{E}(\bar{F})
$$

by:

$$
\alpha_{0}(t)=\frac{a}{b} \quad \text { and } \quad \alpha_{i}(t)=\frac{b}{c_{i}}
$$

for $i=1,2$ or 3 .
From the above description, one sees that the center $Z_{L_{E}}$ of the Levi subgroup $L_{E}$ is such that

$$
X_{*}\left(Z_{L_{E}}\right)=\left\langle\mu_{1}, \mu_{2}, \mu_{3}\right\rangle .
$$

Moreover, the group $\operatorname{Hom}_{F}\left(L_{E}, \mathbb{G}_{m}\right)$ can be described as follows. For $1 \leq i \leq 3$, let us define $\chi_{i} \in X^{*}\left(L_{E}\right)$ by

$$
\chi_{i}\left(g, c_{1}, c_{2}, c_{3}\right)=\frac{\operatorname{det}(g) \cdot c_{i}}{c_{1} c_{2} c_{3}}
$$

Then the elements $\chi_{i}$ form a $\mathbb{Z}$-basis of $X^{*}\left(L_{E}\right)$. On restriction to $T_{E}$, we have:

$$
\chi_{i}=\beta_{0}-\alpha_{0}-\alpha_{i} \quad \text { for } i=1,2 \text { or } 3
$$

Now taking into account the Galois action, we have:

$$
\operatorname{Hom}_{F}\left(L_{E}, \mathbb{G}_{m}\right)=\left\{\begin{array}{l}
\mathbb{Z} \chi_{1} \oplus \mathbb{Z} \chi_{2} \oplus \mathbb{Z} \chi_{3}, \text { if } E=F \times F \times F \\
\mathbb{Z} \chi_{1} \oplus \mathbb{Z} \cdot\left(\chi_{2} \oplus \chi_{3}\right), \text { if } E=F \times K \\
\mathbb{Z} \cdot\left(\chi_{1}+\chi_{2}+\chi_{3}\right), \text { if } E \text { is a field }
\end{array}\right.
$$

Observe that the character

$$
\nu_{L_{E}}=\chi_{1}+\chi_{2}+\chi_{3}
$$

always belong to $\operatorname{Hom}_{F}\left(L_{E}, \mathbb{G}_{m}\right)$. Moreover, the modulus character of $Q_{E}$ is

$$
\delta_{Q_{E}}=3 \cdot \nu_{L_{E}}
$$

2.6. Eisenstein series. With the above description of $\operatorname{Hom}_{F}\left(L_{E}, \mathbb{G}_{m}\right)$, we can now define an unramified character of $L_{E}\left(F_{v}\right)$ or $L_{E}(\mathbb{A})$ :

$$
\chi_{\underline{s}}=\left\{\begin{array}{l}
\left|\chi_{1}\right|^{1+s_{1}} \cdot\left|\chi_{2}\right|^{1+s_{2}} \cdot\left|\chi_{3}\right|^{1+s_{3}}, \text { if } E=F^{3} ; \\
\left|\chi_{1}\right|^{1+s_{1}} \cdot\left|\chi_{2}+\chi_{3}\right|^{1+s_{23}}, \text { if } E=F \times K ; \\
\left|\chi_{1}+\chi_{2}+\chi_{3}\right|^{1+s}, \text { if } E \text { is a field. }
\end{array}\right.
$$

We can then consider the degenerate principal series

$$
I_{Q_{E}}(\underline{s})=I n d_{Q_{E}}^{G_{E}} \chi_{\underline{s}} \quad \text { (unnormalized induction). }
$$

The associated Eisenstein series is given by

$$
E(\Phi, \underline{s}, g)=\sum_{\gamma \in Q_{E}(F) \backslash G_{E}(F)} \Phi_{s}(\gamma g)
$$

for a standard section $\Phi_{s} \in I_{Q_{E}}(\underline{s})$. This sum converges absolutely at $\underline{s}=\left(s_{1}, s_{2}, s_{3}\right)$ if $s_{i}+s_{i+1} \gg s_{i+2}$ for $i=1,2,3$ and where the subscripts are taken modulo 3. We let $\Omega_{R}$ denote the region consisting of those $\underline{s}$ satisfying $s_{i}+s_{i+1}-s_{i+2} \geq R$, for $i=1,2,3$.

This Eisenstein series is one of the ingredients in our Rankin-Selberg integral.
2.7. Geomeric description of $\operatorname{Spin}_{8}(\mathbb{C})$. At this point, we should describe precisely what we mean by the representation $V_{\text {Spin }}$ of $\operatorname{Spin}_{8}(\mathbb{C}) \rtimes S_{3}$. For this, it is necessary to give a geometric description of the split group $\mathrm{Spin}_{8}$. This (very beautiful) description works over an arbitrary field $k$ and can be found, for example, in the paper [Gr].

Let $\mathbb{O}$ be the split octonion algebra over $k$. Then $\mathbb{O}$ comes equipped with

- an anti-involution $x \mapsto \bar{x}$;
- a (quadratic) norm form $N: \mathbb{O} \longrightarrow k$ such that $N(x)=x \cdot \bar{x}=\bar{x} \cdot x$;
- a (linear) trace form $\operatorname{Tr}: \mathbb{O} \longrightarrow k$ such that $\operatorname{Tr}(x)=x+\bar{x}$.

In particular, we may consider the special orthogonal group $S O(\mathbb{O}, N)$.
Though the multiplication in $\mathbb{O}$ is not associative, the symbol $\operatorname{Tr}(x y z)$ is well-defined and satisfies

$$
\operatorname{Tr}(x y z)=\operatorname{Tr}(y z x)=\operatorname{Tr}(z x y) .
$$

Now the group $\mathrm{Spin}_{8}$ over $k$ can be described by:

$$
\operatorname{Spin}_{8}=\left\{\left(g_{1}, g_{2}, g_{3}\right) \in S O(\mathbb{O}, N)^{3}: \operatorname{Tr}\left(\left(g_{1} x\right)\left(g_{2} y\right)\left(g_{3} z\right)\right)=\operatorname{Tr}(x y z) \quad \text { for all } x, y, z \in \mathbb{O}\right\} .
$$

How does $S_{3}$ act on $\operatorname{Spin}_{8}$ ? A naive guess is that it acts on an element $\left(g_{1}, g_{2}, g_{3}\right)$ by permuting its three components. However, this is not the case. By the above, it is true that cyclic permutations of $\left(g_{1}, g_{2}, g_{3}\right)$ do preserve the group $\operatorname{Spin}_{8}$ so that the 3-cycles in $S_{3}$ act in the natural way. A transposition, on the other hand, acts as follows:

$$
\sigma:\left(g_{1}, g_{2}, g_{3}\right) \mapsto\left(\hat{g_{2}}, \hat{g_{1}}, \hat{g_{3}}\right)
$$

where

$$
\hat{g}: x \mapsto \overline{g(\bar{x})}
$$

2.8. The representation $V_{\text {Spin }}$. By the above description of $\operatorname{Spin}_{8}$, one sees that $\operatorname{Spin}_{8}$ acts on $\mathbb{O} \oplus \mathbb{O} \oplus \mathbb{O}$. The action preserves each factor of $\mathbb{O}$ and gives the three irreducible representations of dimension 8 . We can now describe the extension of this action to the semidirect product $\mathrm{Spin}_{8} \rtimes S_{3}$. Not surprisingly, the 3 -cycles in $S_{3}$ acts on $\mathbb{O} \oplus(\mathbb{O} \oplus(\mathbb{O}$ by cyclic permutation. The action of the transposition $\sigma$, however, is given by:

$$
\sigma:(x, y, z) \mapsto(\bar{y}, \bar{x}, \bar{z})
$$

It is easy to show that this does define an action of $\operatorname{Spin}_{8} \rtimes S_{3}$. This representation $V_{\text {Spin }}$ is what we use in the definition of the Spin L-function.

We remark that one could consider the twist of $V_{\text {Spin }}$ by the sign representation of $S_{3}$ and thus obtain a different $L$-function. Though we stated our main theorem in the introduction for $L\left(\underline{s}, \pi, V_{\text {Spin }, E}\right)$, a slight twist of our Rankin-Selberg integral will give an analogous statement for $L\left(\underline{s}, \pi, V_{\text {Spin }, E} \otimes(\operatorname{sign})\right)$.

## 3. Minimal Representation

Now we come to another ingredient in our Rankin-Selberg integral. For each local field $F_{v}$, the group $G_{E}\left(F_{v}\right)$ has a so-called minimal representation $\Pi_{E_{v}}$, first studied by Kazhdan $[K]$. It is the analog of the Weil representation of the metaplectic group.
3.1. Local minimal representation. To describe this minimal representation, let $r: S_{3} \rightarrow$ $\mathrm{GL}_{2}(\mathbb{C})$ denote the 2-dimensional irreducible representation of $S_{3}$. Then the composite $r \circ$ $\rho_{E, v}$ is a 2-dimensional representation of $\operatorname{Gal}\left(\bar{F}_{v} / F_{v}\right)$. By Jacquet-Langlands, this is the Lparameter of an irreducible admissible representation $r_{E, v}$ of $\mathrm{GL}_{2}\left(F_{v}\right)$ given by the following table:

| $E_{v}$ | $r_{E, v}$ |
| :--- | :--- |
| $F_{v}^{3}$ | $\pi(1,1)$ |
| $F_{v} \times K_{v}$ | $\pi\left(1, \chi_{K_{v}}\right)$ |
| Galois field | $\pi\left(\chi_{E_{v}}, \chi_{E_{v}}^{-1}\right)$ |
| non-Galois field | monomial supercuspidal |

Here, $\chi_{K_{v}}$ and $\chi_{E_{v}}$ denote the characters of $F_{v}^{\times}$associated to $K_{v}$ and $E_{v}$ by local class field theory, and $\pi\left(\mu_{1}, \mu_{2}\right)$ denotes the representation of $\mathrm{GL}_{2}\left(F_{v}\right)$ unitarily induced from the character $\mu_{1} \times \mu_{2}$ of the diagonal torus. The central character of $r_{E, v}$ is precisely the quadratic character $\chi_{K_{E, v}}$, and we define a representation of

$$
L_{E}\left(F_{v}\right) \cong\left(\mathrm{GL}_{2}\left(F_{v}\right) \times E_{v}^{\times}\right) / \Delta F_{v}^{\times}
$$

by:

$$
\sigma_{E, v}=r_{E, v} \boxtimes\left(\chi_{K_{E}, v} \circ \mathbb{N}_{E_{v} / F_{v}}\right)
$$

Here and elsewhere, $\boxtimes$ refers to "outer tensor product".
Now $\Pi_{E, v}$ is the unique irreducible submodule of the induced representation

$$
\operatorname{Ind} Q_{Q_{E}}^{G_{E}} \delta_{Q_{E_{v}}}^{1 / 6} \sigma_{E, v} \quad \text { (unnormalized induction). }
$$

If $E_{v}$ is unramified, then $\Pi_{E_{v}}$ is a spherical representation.
3.2. Automorphic realization. Let $\Pi_{E}=\otimes_{v} \Pi_{E_{v}}$ be the global minimal representation of $G_{E}(\mathbb{A})$. It is known (cf. [GGJ]) that there is a unique $G_{E}$-equivariant embedding

$$
\theta: \Pi_{E} \longrightarrow \mathcal{A}\left(G_{E}\right)
$$

which is constructed using residues of Eisenstein series. We need to know some properties of this minimal representation.
3.3. Fourier coefficients. Fix a non-trivial unitary character $\psi_{v}$ of $F_{v}$. Using $\psi_{v}$, the Killing form and the exponential map, the unitary characters of $N_{E}\left(F_{v}\right)$ can be parametrized by elements of $\bar{N}_{E}\left(F_{v}\right) / \bar{Z}\left(F_{v}\right)$. Let $\chi_{v}$ be a non-trivial unitary character of $N_{E}\left(F_{v}\right)$. Then the key property of $\Pi_{E_{v}}$ is that:

$$
\operatorname{dim} \operatorname{Hom}_{N_{E}\left(F_{v}\right)}\left(\Pi_{E_{v}}, \mathbb{C}_{\chi_{v}}\right)=\left\{\begin{array}{l}
1, \text { if } \chi_{v} \in \Omega_{E}\left(F_{v}\right) \\
0, \text { if not. }
\end{array}\right.
$$

Here $\Omega_{E}\left(F_{v}\right)$ is as introduced in (2.4) and parametrizes the minimal non-trivial $M_{E}\left(F_{v}\right)$-orbit of characters of $N_{E}\left(F_{v}\right)$. An element in $\Omega_{E}\left(F_{v}\right)$ is the character $\chi_{0}$ such that

$$
\left.\chi_{0}\right|_{U_{\beta_{0}-\alpha_{0}}}=\psi_{v}
$$

and $\left.\chi_{0}\right|_{U_{\gamma}}$ is trivial for any other root subgroup $U_{\gamma} \subset N_{E}$. In terms of the description of $\Omega_{E}$ given in (2.4),

$$
\chi_{0}=(0,0,0,1)
$$

These local results on $\operatorname{Hom}_{N_{E}\left(F_{v}\right)}\left(\Pi_{E_{v}}, \mathbb{C}_{\chi_{v}}\right)$ imply:
Proposition 3.3.1. (i) For $f \in \Pi_{E}$, the Fourier expansion of $\theta(f)$ along $N_{E}$ is supported on the (Zariski) closure of the minimal non-trivial $M_{E}(F)$-orbit $\Omega_{E}(F)$ of unitary characters of $N_{E}(\mathbb{A})$ trivial on $N_{E}(F)$.
(ii) For $\chi \in \Omega_{E}(F)$, the Fourier coefficient $\theta(f)_{N_{E}, \chi}$ has an Euler product expansion in terms of the local functionals. Namely, if $f=\otimes_{v} f_{v}$, then

$$
\theta(f)_{N_{E}, \chi}(g)=\prod_{v} L_{\chi_{v}}\left(g_{v} \cdot f_{v}\right)
$$

where each $L_{\chi_{v}}$ is a non-zero element of the 1-dimensional space $\operatorname{Hom}_{N_{E}\left(F_{v}\right)}\left(\Pi_{v}, \mathbb{C}_{\chi_{v}}\right)$, normalized for almost all $v$ by the requirement that $L_{\chi_{v}}\left(f_{v}^{0}\right)=1$ if $f_{v}^{0}$ is the normalized spherical vector in $\Pi_{E, v}$.
3.4. An explicit formula. Finally, we come to the only new result in this section. We have fixed a Chevalley-Steinberg system of épinglage for $G_{E}$ over $F$. For almost all finite places $v$ of $F$, this system of épinglage gives $G_{E}$ an $\mathcal{O}_{F_{v}}$-structure such that $G_{E}\left(\mathcal{O}_{F_{v}}\right)$ is a hyperspecial maximal compact subgroup of $G_{E}\left(F_{v}\right)$.

Fix such a non-archimedean place $v$, and let $L_{\chi_{0}}$ be a non-zero element of the 1-dimensional space $\operatorname{Hom}_{N_{E}}\left(\Pi_{E_{v}}, \mathbb{C}_{\chi_{0}}\right)$. We have the following explicit formula:
Proposition 3.4.1. Assume that $v$ is a finite place such that $G_{E}\left(\mathcal{O}_{F_{v}}\right)$ is hyperspecial (so that $E_{v}$ is unramified) and let $f_{0}$ be a non-zero $G_{E}\left(\mathcal{O}_{F_{v}}\right)$-spherical vector in $\Pi_{E_{v}}$. For $t \in T_{E}\left(F_{v}\right)$, let us set

$$
k_{0}(t)=\operatorname{ord}\left(\beta_{0}(t) / \alpha_{0}(t)\right)
$$

If $\chi_{0}$ has conductor zero (with respect to the $\mathcal{O}_{F_{v}}$-structure on $N_{E}\left(F_{v}\right)$ ), then

$$
L_{\chi_{0}}\left(t f_{0}\right)=0 \quad \text { if } k_{0}(t)<0
$$

However, $L_{\chi_{0}}$ is non-zero on $f_{0}$, so that we may normalize $L_{\chi_{0}}$ by setting $L_{\chi_{0}}\left(f_{0}\right)=1$. Then, if $k_{0}(t) \geq 0$, we have:
(i) If $E_{v}=F_{v}^{3}$, then

$$
L_{\chi_{0}}\left(t \cdot f_{0}\right)=\left|\beta_{0}(t)\right| \cdot\left(k_{0}(t)+1\right) .
$$

(ii) If $E_{v}=F_{v} \times K_{v}$, then

$$
L_{\chi_{0}}\left(t \cdot f_{0}\right)=\chi_{K_{v}}\left(\beta_{0}(t)\right) \cdot\left|\beta_{0}(t)\right| \cdot \epsilon_{2}\left(k_{0}(t)+1\right)
$$

where $\epsilon_{2}(\cdot)$ is the Legendre symbol $(\dot{\overline{2}})$ (with $\epsilon_{2}(0)=0$ ).
(iii) If $E$ is the unramified cubic field extension of $F_{v}$, then

$$
L_{\chi_{0}}\left(t \cdot f_{0}\right)=\left|\beta_{0}(t)\right| \cdot \epsilon_{3}\left(k_{0}(t)+1\right),
$$

where $\epsilon_{3}(\cdot)$ is the Legendre symbol $(\dot{\overline{3}}) \quad$ (with $\epsilon_{3}(0)=0$ ).
Proof. In the split case, this is a result of Kazhdan-Polishchuk [KP]. We give a different proof which covers the non-split case as well.

Choose a representative in $G_{E}\left(\mathcal{O}_{F_{v}}\right)$ of the Weyl group element $w_{\alpha}$; we denote this representative by $w_{\alpha}$ also. Then $w_{\alpha}$ normalizes $N_{E}$ and it will be more convenient to replace $\chi_{0}$ and $L_{\chi_{0}}$ by $\chi_{0}^{\prime}=w_{\alpha} \chi_{0}$ and $L^{\prime}=w_{\alpha} L_{\chi_{0}}$. Let us give an explicit construction of $L^{\prime}$.

By definition, the linear functional $L^{\prime}$ factors through the quotient $\left(\Pi_{E_{v}}\right)_{N_{E}, \chi_{0}^{\prime}}$ which is 1-dimensional. The root subgroup $U_{\alpha}\left(F_{v}\right)$ normalizes $N_{E}\left(F_{v}\right)$ and fixes the character $\chi_{0}^{\prime}$. It thus acts on $\left(\Pi_{E_{v}}\right)_{N_{E}, \chi_{0}^{\prime}}$ and this action is in fact trivial. Thus, we see that $L^{\prime}$ factors as:

$$
\Pi_{E_{v}} \longrightarrow\left(\Pi_{E_{v}}\right)_{U_{E}} \xrightarrow{l_{\psi_{v}}} \mathbb{C}
$$

where

$$
l_{\psi_{v}} \in \operatorname{Hom}_{U_{\beta}}\left(\left(\Pi_{E_{v}}\right)_{U_{E}}, \mathbb{C}_{\psi_{v}}\right) .
$$

Now we recall that $\Pi_{E, v}$ is the unique irreducible submodule of

$$
\operatorname{Ind} d_{Q_{E}}^{G_{E}} \delta_{Q_{E_{v}}}^{1 / 6} \sigma_{E, v}
$$

By Frobenius reciprocity, there is a natural $Q_{E_{v}}$-equivariant map

$$
p: \Pi_{E_{v}} \longrightarrow \delta_{Q_{E_{v}}}^{1 / 6} \otimes \sigma_{E, v}
$$

which is simply given

$$
p(f)=f(1) \quad \text { for } f \in \Pi_{E_{v}} .
$$

Thus $L^{\prime}$ is the composite of $p$ with the unique Whittaker functional $l_{\psi_{v}}$ of $\sigma_{E, v}$.
Now take a torus element

$$
t=\left(\left(\begin{array}{ll}
a & 0 \\
0 & b
\end{array}\right), x\right) \in T_{E}\left(F_{v}\right) \subset L_{E}\left(F_{v}\right)=\left(G L_{2}\left(F_{v}\right) \times E_{v}^{\times}\right) / \Delta F_{v}^{\times} .
$$

Then

$$
\begin{aligned}
& L^{\prime}\left(t \cdot f_{0}\right)=l_{\psi_{v}}\left(f_{0}(t)\right)=\delta_{Q_{E_{v}}}(t)^{1 / 6} \cdot l_{\psi_{v}}\left(\sigma_{E, v}(t)\left(f_{0}(1)\right)\right) \\
&=\frac{|a b|^{3 / 2}}{|\mathbb{N}(x)|} \cdot \chi_{K_{E}}(\mathbb{N}(x)) \cdot W_{\psi_{v}}\left(\begin{array}{cc}
a & 0 \\
0 & b
\end{array}\right) \\
& 11
\end{aligned}
$$

where $W_{\psi_{v}}$ is the normalized spherical Whittaker function of the representation $r_{E, v}$ of $G L_{2}\left(F_{v}\right)$. The formula for $W_{\psi_{v}}$ on a torus element is well-known:

$$
|a / b|^{-1 / 2} \cdot W_{\psi_{v}}\left(\begin{array}{cc}
a & 0 \\
0 & b
\end{array}\right)=\chi_{K_{E}}(b) \cdot \epsilon(\operatorname{ord}(a / b))
$$

where

$$
\epsilon(k):=\left\{\begin{array}{l}
k+1, \text { if } E_{v}=F_{v}^{3} \\
\epsilon_{2}(k+1), \text { if } E_{v}=F_{v} \times K_{v} \\
\epsilon_{3}(k+1), \text { if } E_{v} \text { is the unramified cubic field. }
\end{array}\right.
$$

Thus, we deduce that

$$
L^{\prime}\left(t \cdot f_{0}\right)=\frac{\left|a^{2} b\right|}{|\mathbb{N}(x)|} \cdot \chi_{K_{E}}(b \cdot \mathbb{N}(x)) \cdot \epsilon(\operatorname{ord}(a / b))=\left|\beta_{0}(t)\right| \cdot \chi_{K_{E}}\left(\beta_{0}(t)\right) \cdot \epsilon\left(\operatorname{ord}\left(\alpha_{0}(t)\right)\right.
$$

Conjugating back by $w_{\alpha}$ gives the desired result, since $w_{\alpha}$ fixes $\beta_{0}$ and sends $\alpha_{0}$ to $\beta_{0}-\alpha_{0}$. The proposition is proved.

For a general vector $f \in \Pi_{E}$, the above construction of the linear functional $L_{\chi_{0}}$ allows one to obtain the asymptotics of the function $(t, k) \mapsto L_{\chi_{0}}(t k \cdot f)$ on $T_{E}\left(F_{v}\right) \times K_{v}$. This is needed in Section 9; see Lemma 9.2.2. Further, in Lemma 9.5.1, we shall give another description of the functional $L_{\chi_{0}}$ by means of a local Fourier-Jacobi map.

## 4. A Rankin-Selberg Integral

We are now ready to write down our Rankin-Selberg integral. Let $\pi$ be a globally generic cuspidal representation of $G_{E}(\mathbb{A})$. For $\varphi \in \pi, \Phi_{s} \in I_{Q_{E}}(\underline{s})$ and $f \in \Pi_{E}$, we set

$$
Z_{E}(\varphi, \Phi, f, \underline{s})=\int_{G_{E}(F) \backslash G_{E}(\mathbb{A})} \varphi(g) \cdot \overline{\theta(f)(g)} \cdot E(\Phi, \underline{s}, g) d g
$$

The purpose of this section is to unfold this Rankin-Selberg integral of Shimura type.
Theorem 4.0.2. Let $V_{E}^{\prime}$ be the maximal unipotent subgroup of $G_{E}$ given by

$$
V_{E}^{\prime}=w_{\alpha} V_{E} w_{\alpha}^{-1}=N_{E} \rtimes U_{-\alpha} .
$$

It is the unipotent radical of the Borel subgroup $B_{E}^{\prime}=w_{\alpha} B_{E} w_{\alpha}^{-1}$, with associated simple roots $-\alpha_{i}(i=1,2,3)$ and $\beta_{0}-\alpha_{0}$. Let $\psi$ be the generic unitary character of $V_{E}^{\prime}(\mathbb{A})$ which is non-trivial on the associated relative simple root subgroups. Then we have:

$$
Z_{E}(\varphi, \Phi, f, \underline{s})=\int_{V_{E}^{\prime}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \varphi_{V_{E}^{\prime}, \psi}(g) \cdot \overline{\theta(f)_{N_{E}, \chi_{0}}(g)} \cdot \widehat{\Phi_{s}}(g) d g
$$

where

$$
\widehat{\Phi_{s}}(g)=\int_{U_{-\alpha}(\mathbb{A})} \Phi_{s}(u g) \cdot \psi(u) d u
$$

Proof. Unfolding the Eisenstein series in the range of absolute convergence, we get

$$
\begin{aligned}
& Z_{E}(\varphi, \Phi, f, \underline{s}) \\
= & \int_{Q_{E}(F) \backslash G_{E}(\mathbb{A})} \varphi(g) \cdot \overline{\theta(f)(g)} \cdot \Phi_{s}(g) d g \\
= & \int_{Q_{E}(F) U_{E}^{(1)}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \Phi_{s}(g) \cdot\left(\int_{U_{E}^{(1)}(F) \backslash U_{E}^{(1)}(\mathbb{A})} \varphi(z g) \cdot \overline{\theta(f)(z g)} d z\right) d g \\
= & \int_{Q_{E}(F) U_{E}^{(1)}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \Phi_{s}(g) \cdot \varphi_{U_{E}^{(1)}}(g) \cdot \overline{\theta(f)_{U_{E}^{(1)}(g)}} d g \\
& +\int_{Q_{E}(F) U_{E}^{(1)}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \Phi_{s}(g) \cdot\left(\sum_{\chi \neq 1} \varphi_{U_{E}^{(1)}, \chi}(g) \cdot \overline{\theta(f)_{U_{E}^{(1)}, \chi}(g)}\right) d g .
\end{aligned}
$$

Lemma 4.0.3. The first term in the last sum is zero.
Proof. The argument is easiest when $E$ is a field and is more complicated in the split case; it ultimately relies on the fact that $\varphi$ is cuspidal.

We shall only give the argument when $E$ is a field. Consider the Fourier expansion of $\theta(f)_{U_{E}^{(1)}}$ along $N_{E}(F) \backslash N_{E}(\mathbb{A})$. Because $\theta(f)_{U_{E}^{(1)}}$ is left-invariant under $U_{\beta_{0}-\alpha_{0}}(\mathbb{A})$, this has the form

$$
\theta(f)_{U_{E}^{(1)}}=\sum_{\chi \in \Omega_{E}(F):\left.\chi\right|_{U_{\beta_{0}-\alpha_{0}}}=1} \theta(f)_{N_{E}, \chi} .
$$

In other words, the character $\chi$ intervening in the above sum is represented by $(a, x, y, 0) \in$ $\Omega_{E}(F)$. But if $E$ is a field, then the only elements of this form in $\Omega_{E}(F)$ are ( $a, 0,0,0$ ) with $a \in F$.

Consider first the constant term $\theta(f)_{N_{E}}$ in the Fourier expansion, which corresponds to $a=0$. By [GGJ, Prop. 5.3(iv)], the restriction of $\theta(f)_{N_{E}}$ to $M_{E}$ lies in the span of two automorphic characters of $M_{E}(\mathbb{A})$. Since this is the case for any $f \in \Pi_{E}$, we see that $\theta(f)_{N_{E}}$ is left-invariant under $U_{\alpha}(\mathbb{A})$, and thus under $U_{E}(\mathbb{A})$.

Now consider the other terms in the Fourier expansion of $\theta(f)_{U_{E}^{(1)}}$. If $\psi_{a}$ is the character associated to $(a, 0,0,0) \in \Omega_{E}(F)$ with $a \neq 0$, then $U_{\alpha}(\mathbb{A})$ normalizes $N_{E}(\mathbb{A})$ and fixes the character $\psi_{a}$. Thus the Fourier coefficient $\theta(f)_{N_{E}, \psi_{a}}$ is left-invariant under $U_{\alpha}(\mathbb{A})$ (we have used the local analog of this fact in the proof of Prop. 3.4.1). Hence $\theta(f)_{N_{E}, \psi_{a}}$ is also left-invariant under $U_{E}(\mathbb{A})$. Hence we have:

$$
\begin{aligned}
& \int_{Q_{E}(F) U_{E}^{(1)}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \Phi_{s}(g) \cdot \varphi_{U_{E}^{(1)}}(g) \cdot\left(\sum_{a \in F} \overline{\theta(f)_{N_{E}, \psi_{a}}(g)}\right) d g \\
= & \int_{L_{E}(F) U_{E}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \Phi_{s}(g) \cdot\left(\sum_{a \in F} \overline{\theta(f)_{N_{E}, \psi_{a}}(g)}\right) \cdot\left(\int_{U_{E}(F) \backslash U_{E}(\mathbb{A})} \varphi(u g) d u\right) d g \\
= & 0
\end{aligned}
$$

since $\varphi$ is cuspidal. This proves the lemma when $E$ is a field.

When $E$ is not a field, there will be more terms intervening in the Fourier expansion of $\theta(f)_{U_{E}^{(1)}}$. Thankfully, there are also more $F$-rational standard parabolics in $G_{E}$. These other terms in the Fourier expansion ultimately lead to the constant terms of $\varphi$ along these other parabolics. We omit the details.

On the other hand, the non-trivial characters of $U_{E}^{(1)}(\mathbb{A})$ in the second term are permuted transitively by $Q_{E}(F)$. If we let $\chi_{0}$ denote the character which is trivial on $U_{\beta_{0}}$ and nontrivial on $U_{\beta_{0}-\alpha_{0}}$, then the stabilizer of $\chi_{0}$ in $Q_{E}(F)$ is the subgroup $A_{E} V_{E}$ where $A_{E}$ is the 3-dimensional torus in $T_{E}$ such that

$$
X_{*}\left(A_{E}\right)=\left\langle\mu_{1}-\mu_{0}, \mu_{2}-\mu_{0}, \mu_{3}-\mu_{0}\right\rangle
$$

In view of this, we have:

$$
\begin{aligned}
& Z_{E}(\varphi, \Phi, f, \underline{s}) \\
= & \int_{Q_{E}(F) U_{E}^{(1)}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \Phi_{s}(g) \cdot\left(\sum_{\gamma \in A_{E}(F) V_{E}(F) \backslash Q_{E}(F)} \varphi_{U_{E}^{(1)}, \chi_{0}}(\gamma g) \cdot \overline{\theta(f)_{U_{E}^{(1)}, \chi_{0}}(\gamma g)}\right) d g \\
= & \int_{A_{E}(F) V_{E}(F) U_{E}^{(1)}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \Phi_{s}(g) \cdot \varphi_{U_{E}^{(1)}, \chi_{0}}(g) \cdot \overline{\theta(f)_{U_{E}^{(1)}, \chi_{0}}(g)} d g \\
= & \int_{A_{E}(F) U_{\alpha}(F) N_{E}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \Phi_{s}(g) \cdot\left(\int_{N_{E}(F) U_{E}^{(1)}(\mathbb{A}) \backslash N_{E}(\mathbb{A})} \varphi_{U_{E}^{(1)}, \chi_{0}}(n g) \cdot \overline{\theta(f)_{U_{E}^{(1)}, \chi_{0}}(n g)}\right) d g
\end{aligned}
$$

By Prop. 3.3.1(i), we see that

$$
\theta(f)_{U_{E}^{(1)}, \chi_{0}}(g)=\sum_{x \in E} \theta(f)_{N_{E}, \psi_{x}}(g) .
$$

Here $\psi_{x}$ corresponds to the element

$$
\left(\mathbb{N}(x), x^{\#}, x, 1\right) \in \Omega_{E}(F)
$$

Moreover, the group $U_{\alpha}(F) \cong E$ acts simply transitively on the set $\left\{\psi_{x}: x \in E\right\}$. Note that if $x=0$, then $\psi_{0}$ is simply the trivial extension of $\chi_{0}$ from $U_{E}^{(1)}$ to $N_{E}$. Thus,

$$
\begin{aligned}
Z_{E}(\varphi, \Phi, f, \underline{s}) & =\int_{A_{E}(F) U_{\alpha}(F) N_{E}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \Phi_{s}(g) \cdot\left(\sum_{x \in E} \varphi_{N_{E}, \psi_{x}}(g) \cdot \overline{\theta(f)_{N_{E}, \psi_{x}}(g)}\right) d g \\
& =\int_{A_{E}(F) N_{E}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \Phi_{s}(g) \cdot \varphi_{N_{E}, \psi_{0}}(g) \cdot \overline{\theta(f)_{N_{E}, \psi_{0}}(g)} d g
\end{aligned}
$$

Now consider the Fourier expansion of $\varphi_{N_{E}, \psi_{0}}$ along $U_{-\alpha}$. Because $\varphi$ is cuspidal, this takes the form:

$$
\varphi_{N_{E}, \psi_{0}}(g)=\sum_{\gamma \in A_{E}(F)} \varphi_{V_{E}^{\prime}, \psi}(\gamma g),
$$

where $\psi$ is a generic character of $V_{E}^{\prime}$. Substituting this into the last expression for our Rankin-Selberg integral, and using the fact that $\theta(f)_{N_{E}, \psi_{0}}$ is left-invariant under $U_{-\alpha}(\mathbb{A})$, we
obtain:

$$
\begin{aligned}
& Z_{E}(\varphi, \Phi, f, \underline{s}) \\
= & \int_{A_{E}(F) N_{E}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \Phi_{s}(g) \cdot \overline{\theta(f)_{N_{E}, \psi_{0}}(g)} \cdot\left(\sum_{\gamma \in A_{E}(F)} \varphi_{V_{E}^{\prime}, \psi}(\gamma g)\right) d g \\
= & \int_{N_{E}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \Phi_{s}(g) \cdot \overline{\theta(f)_{N_{E}, \psi_{0}}(g)} \cdot \varphi_{V_{E}^{\prime}, \psi}(g) d g \\
= & \int_{V_{E}^{\prime}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \overline{\theta(f)_{N_{E}, \psi_{0}}(g)} \cdot\left(\int_{U_{-\alpha}(\mathbb{A})} \varphi_{V_{E}^{\prime}, \psi}(u g) \cdot \Phi_{s}(u g) d u\right) d g \\
= & \int_{V_{E}^{\prime}(\mathbb{A}) \backslash G_{E}(\mathbb{A})} \overline{\theta(f)_{N_{E}, \psi_{0}}(g)} \cdot \varphi_{V_{E}^{\prime}, \psi}(g) \cdot \widehat{\Phi_{s}}(g) d g .
\end{aligned}
$$

Theorem 4.0.2 is proved.

## 5. Local Zeta Integral

After Theorem 4.0.2, we see that

$$
Z_{E}(\varphi, \Phi, f, \underline{s})=\prod_{v} Z_{v}\left(\varphi_{v}, \Phi_{v}, f_{v}, \underline{s}\right)
$$

where

$$
Z_{v}\left(\varphi_{v}, \Phi_{v}, f_{v}, \underline{s}\right)=\int_{V_{E}^{\prime}\left(F_{v}\right) \backslash G_{E}\left(F_{v}\right)} W_{V_{E}^{\prime}, \psi}\left(g \cdot \varphi_{v}\right) \cdot \overline{L_{\chi_{0}}\left(g \cdot f_{v}\right)} \cdot \widehat{\Phi_{s, v}}(g) d g
$$

This integral converges when $\underline{s} \in \Omega_{R}$ (cf. (2.6)) for $R \gg 0$, as we shall show in Prop. 9.1.3. The rest of the paper is devoted to the study of this local zeta integral. In particular, we shall compute the local zeta integral explicitly when all the data involved are unramified. The purpose of this section is to provide explicit formulas for the 3 functions appearing in the local integral in the unramified setting.

Assume henceforth that all the data involved in $Z_{v}$ are unramified. Then by the Iwasawa decomposition, we have

$$
Z_{v}\left(\varphi_{v}, \Phi_{v}, f_{v}, s\right)=\int_{S_{E_{v}}\left(F_{v}\right)} W_{V_{E}^{\prime}, \psi}\left(t \cdot \varphi_{v}\right) \cdot \overline{L_{\chi_{0}}\left(t \cdot f_{v}\right)} \cdot \widehat{\Phi_{s, v}}(t) \cdot \delta_{B_{E}^{\prime}}(t)^{-1} d t
$$

where $S_{E_{v}} \subset T_{E_{v}}$ is the maximal $F_{v}$-split torus.
Henceforth, since the setting is entirely local, we shall suppress $v$ from the notations.
5.1. A change of system of simple roots. Because we are looking at the Whittaker functional relative to the Borel subgroup $B_{E}^{\prime}$, it is useful to use the corresponding system of simple roots:

$$
\alpha_{0}^{\prime}=\beta_{0}-\alpha_{0}, \quad \alpha_{i}^{\prime}=-\alpha_{i}, \quad \text { for } i=1,2 \text { and } 3
$$

For this new system of simple roots, the highest root is still

$$
\beta_{0}=2 \alpha_{0}^{\prime}+\alpha_{1}^{\prime}+\alpha_{2}^{\prime}+\alpha_{3}^{\prime}
$$

and the modulus character of $B_{E}^{\prime}$ is given by

$$
\delta_{B_{E}^{\prime}}=\left|6\left(\alpha_{1}^{\prime}+\alpha_{2}^{\prime}+\alpha_{3}^{\prime}\right)+10 \alpha_{0}^{\prime}\right|
$$

The associated fundamental coweights are given by:

$$
\mu_{0}^{\prime}=\mu_{0} \quad \mu_{i}^{\prime}=\mu_{0}-\mu_{i} \quad \text { for } i=1,2 \text { and } 3
$$

As above, let $S_{E} \subset T_{E}$ denote the maximal $F$-split torus. We consider the 3 different cases:

- if $E=F \times F \times F$, then $S_{E}=T_{E}$. An element of $S_{E}(F)$ is of the form

$$
t=\mu_{0}^{\prime}\left(t_{0}\right) \cdot \mu_{1}^{\prime}\left(t_{1}\right) \cdot \mu_{2}^{\prime}\left(t_{2}\right) \cdot \mu_{3}^{\prime}\left(t_{3}\right), \quad \text { with } t_{i} \in F^{\times}
$$

We shall sometimes write this element as $t=\left(t_{0}, t_{1}, t_{2}, t_{3}\right)$.

- if $E=F \times K$ with $K$ a quadratic field, then

$$
X_{*}\left(S_{E}\right)=\left\langle\mu_{0}^{\prime}, \mu_{1}^{\prime}, \mu_{2}^{\prime}+\mu_{3}^{\prime}\right\rangle
$$

An element of $S_{E}(F)$ is of the form

$$
t=\mu_{0}^{\prime}\left(t_{0}\right) \cdot \mu_{1}^{\prime}\left(t_{1}\right) \cdot\left(\mu_{2}^{\prime} \mu_{3}^{\prime}\right)\left(t_{23}\right)
$$

We shall write this element as $t=\left(t_{0}, t_{1}, t_{23}\right)$.

- if $E$ is a cubic field, then

$$
X_{*}\left(S_{E}\right)=\left\langle\mu_{0}^{\prime}, \mu_{1}^{\prime}+\mu_{2}^{\prime}+\mu_{3}^{\prime}\right\rangle
$$

An element of $S_{E}(F)$ is of the form

$$
t=\mu_{0}^{\prime}\left(t_{0}\right) \cdot\left(\mu_{1}^{\prime} \mu_{2}^{\prime} \mu_{3}^{\prime}\right)\left(t_{123}\right)
$$

We shall write this element as $t=\left(t_{0}, t_{123}\right)$.
5.2. Casselman-Shalika formula. For the value of the Whittaker functional $W_{V_{E}^{\prime}, \psi}$ on $S_{E}(F)$, one has the well-known Casselman-Shalika formula. Let us state this precisely since it is slightly more subtle when the group $G_{E}$ is not split. To the best of our knowledge, the first discussion about the interpretation of the Casselman-Shalika formula in terms of an appropriate dual group in non-split cases is due to Tamir [T]. Our description below is somewhat cleaner than that in $[\mathrm{T}]$, since the relative root system involved here is reduced. The description in $[T]$ works for all relative root systems; indeed the case needed in $[\mathrm{T}]$ is that of the quasi-split unitary groups whose relative root system is of type $B C_{n}$.

For $t \in S_{E}(F), W_{V_{E}^{\prime}, \psi}(t \cdot \varphi)$ is zero unless $\left|t_{i}\right| \leq 1$ for all entries $t_{i}$ of $t$, in which case it depends only on the valuation of the $t_{i}$ 's. If we write $t_{i}=\varpi^{k_{i}}$ with $k_{i} \geq 0$, then we denote the corresponding torus element by

$$
t=t(\underline{k})=t\left(k_{0}, k_{1}, k_{2}, k_{3}\right) \text { or } t\left(k_{0}, k_{1}, k_{23}\right) \text { or } t\left(k_{0}, k\right)
$$

in the three respective cases.
Now we need to examine the notion of Satake parameter for a general quasi-split group with reduced relative root system. The Satake parameter $t_{\pi}$ which was used in the definition
 group ${ }^{L} G_{E}$, where $F_{u r}$ is the maximal unramified extension of $F$ in $\bar{F}$. In fact, $t_{\pi}$ lies in the coset $G_{E}^{\vee}$. Frob of the Frobenius element Frob. Thus we may write

$$
t_{\pi}=\left(s_{\pi}, \operatorname{Frob}\right) \in G_{E}^{\vee} \rtimes \operatorname{Gal}\left(F_{u r} / F\right)
$$

where now $s_{\pi}$ is well-defined up to Frob-conjugacy as an element of $G_{E}^{\vee}$. Under Frobconjugation, we may assume further that

$$
s_{\pi} \in H_{E}^{\prime}:=\left(G_{E}^{\vee}\right)^{\operatorname{Gal}\left(F_{u r} / F\right)},
$$

in which case $s_{\pi}$ is well-defined up to conjugacy in $H_{E}^{\prime}$. Note that

$$
H_{E}^{\prime}=\left\{\begin{array}{l}
\operatorname{Spin}_{8}(\mathbb{C}), \text { if } E=F^{3} ; \\
\operatorname{Spin}_{7}(\mathbb{C}), \text { if } E=F \times K \\
G_{2}(\mathbb{C}), \text { if } E \text { is a field. }
\end{array}\right.
$$

However, the Casselman-Shalika formula is not interpreted in terms of the element $t_{\pi}$ or the element $s_{\pi}$. Rather, the group $G_{E}$ contains a connected $F$-split subgroup $H_{E}$ whose root system is equal to the relative root system of $G_{E}$ and whose maximal split torus is the maximal $F$-split subtorus $S_{E}$ of $T_{E}$. In our case,

$$
H_{E}=\left\{\begin{array}{l}
G_{E}, \text { if } E=F^{3} \\
S O_{7}, \text { if } E=F \times K \\
G_{2}, \text { if } E \text { is a field }
\end{array}\right.
$$

so that its dual group is

$$
H_{E}^{\vee}=\left\{\begin{array}{l}
\operatorname{Spin}_{8}(\mathbb{C}), \text { if } E=F^{3} \\
S_{6}(\mathbb{C}), \text { if } E=F \times K \\
G_{2}(\mathbb{C}), \text { if } E \text { is a field }
\end{array}\right.
$$

From the diagram of inclusions

we obtain on the dual side:

$$
\begin{aligned}
& T_{E}^{\vee} \longrightarrow G_{E}^{\vee} \\
& \iota^{*} \mid \\
& S_{E}^{\vee} \longrightarrow H_{E}^{\vee}
\end{aligned}
$$

where the horizontal arrows are inclusions but $\iota^{*}$ is surjective. Indeed, $\iota^{*}$ restricts to an isogeny

$$
\left(T_{E}^{\vee}\right)^{\operatorname{Gal}\left(F_{u r} / F\right)} \longrightarrow S_{E}^{\vee}
$$

As explained by Borel in [Bo, §6], the map $\iota^{*}$ induces a bijection
\{semisimple Frob-conjugacy classes in $\left.G_{E}^{\vee}\right\} \leftrightarrow\left\{\right.$ semisimple conjugacy classes in $\left.H_{E}^{\vee}\right\}$
We set

$$
\overline{s_{\pi}}=\iota^{*}\left(s_{\pi}\right)
$$

It is the element $\overline{s_{\pi}}$ which intervenes in the Casselman-Shalika formula.
More precisely, the element $t=t(\underline{k})$ corresponds to an element of $X_{*}\left(S_{E}\right) \cong X^{*}\left(S_{E}^{\vee}\right)$ in the dominant chamber and thus gives rise to an irreducible representation $V_{\underline{k}}$ of $H_{E}^{\vee}$. The Casselman-Shalika formula says:

Proposition 5.2.1. With the notations introduced above,

$$
\delta_{B_{E}^{\prime}}(t(\underline{k}))^{-1 / 2} \cdot W_{V_{E}^{\prime}, \psi}(t(\underline{k}) \cdot \varphi)=\operatorname{Tr}\left(\overline{s_{\pi}} \mid V_{\underline{k}}\right) .
$$

In our unramified computations in the following sections, we shall state more precisely what the representation $V_{\underline{k}}$ is for each given $\underline{k}$.
5.3. The map $\iota^{*}$. For the purpose of calculation, we need to understand the map $\iota^{*}$ more explicitly. When $E=F^{3}$, there is nothing to do, since $\iota$ is the identity map and $\overline{s_{\pi}}=s_{\pi}$. We examine the other two cases in turn. Since

$$
X_{*}\left(T_{E}^{\vee}\right)=X^{*}\left(T_{E}\right)=\bigoplus_{i=0}^{3} \mathbb{Z} \alpha_{i}^{\prime},
$$

the element $s_{\pi} \in T_{E}^{\vee}$ is of the form

$$
s_{\pi}=\prod_{i=0}^{3} \alpha_{i}^{\prime}\left(t_{i}\right)
$$

Indeed, if $s_{\pi}$ is assumed to lie in $\left(T_{E}^{\vee}\right)^{\operatorname{Gal}\left(F_{u r} / F\right)}$, then we further have $t_{2}=t_{3}$ (respectively, $t_{1}=t_{2}=t_{3}$ ) when $E=F \times K$ (respectively when $E$ is a field). On the other hand,

$$
X_{*}\left(S_{E}^{\vee}\right)=X^{*}\left(S_{E}\right) \cong\left\{\begin{array}{l}
X^{*}\left(T_{E}\right) /\left\langle\alpha_{2}^{\prime}-\alpha_{3}^{\prime}\right\rangle, \text { if } E=F \times K \\
X^{*}\left(T_{E}\right) /\left\langle\alpha_{1}^{\prime}-\alpha_{2}^{\prime}, \alpha_{2}^{\prime}-\alpha_{3}^{\prime}\right\rangle, \text { if } E \text { is a field. }
\end{array}\right.
$$

If we let $\bar{\alpha}_{i}^{\prime}$ denote the image of $\alpha_{i}^{\prime}$ in the quotient lattice above, then a basis of $X_{*}\left(S_{E}^{\vee}\right)$ is

$$
\left\{\begin{array}{l}
\left\{\bar{\alpha}_{0}^{\prime}, \bar{\alpha}_{1}^{\prime}, \bar{\alpha}_{2}^{\prime}\right\}, \text { if } E=F \times K ; \\
\left\{\bar{\alpha}_{0}^{\prime}, \bar{\alpha}_{1}^{\prime}\right\}, \text { if } E \text { is a field. }
\end{array}\right.
$$

Now we have:

- when $E=F \times K$,

$$
\iota^{*}: \alpha_{0}\left(t_{0}\right) \cdot \alpha_{1}^{\prime}\left(t_{1}\right) \cdot \alpha_{2}^{\prime}\left(t_{2}\right) \cdot \alpha_{3}^{\prime}\left(t_{3}\right) \mapsto \bar{\alpha}_{0}^{\prime}\left(t_{0}\right) \cdot \bar{\alpha}_{1}^{\prime}\left(t_{1}\right) \cdot \bar{\alpha}_{2}^{\prime}\left(t_{2} t_{3}\right)
$$

- when $E$ is a field,

$$
\iota^{*}: \alpha_{0}\left(t_{0}\right) \cdot \alpha_{1}^{\prime}\left(t_{1}\right) \cdot \alpha_{2}^{\prime}\left(t_{2}\right) \cdot \alpha_{3}^{\prime}\left(t_{3}\right) \mapsto \bar{\alpha}_{0}^{\prime}\left(t_{0}\right) \cdot \bar{\alpha}_{1}^{\prime}\left(t_{1} t_{2} t_{3}\right)
$$

5.4. Formula for $L_{\chi_{0}}(t \cdot f)$. On the other hand, Prop. 3.4.1 gives an explicit formula for $L_{\chi_{0}}(t \cdot f)$. For ease of reference, we restate it here:

- when $E=F^{3}$,

$$
L_{\chi_{0}}\left(t\left(k_{0}, k_{1}, k_{2}, k_{3}\right) \cdot f\right)=\left|\beta_{0}(t)\right| \cdot\left(\operatorname{ord}\left(\alpha_{0}^{\prime}(t)\right)+1\right)=q^{-\left(2 k_{0}+k_{1}+k_{2}+k_{3}\right)} \cdot\left(k_{0}+1\right)
$$

- when $E=F \times K$,

$$
L_{\chi_{0}}\left(t\left(k_{0}, k_{1}, k_{23}\right) \cdot f\right)=(-1)^{k_{1}} \cdot q^{-\left(2 k_{0}+k_{1}+2 k_{23}\right)} \cdot \epsilon_{2}\left(k_{0}+1\right) .
$$

- when $E$ is a field,

$$
L_{\chi_{0}}\left(t\left(k_{0}, k_{123}\right) \cdot f\right)=q^{-\left(2 k_{0}+3 k_{123}\right)} \cdot \epsilon_{3}\left(k_{0}+1\right)
$$

5.5. Formula for $\widehat{\Phi_{s}}(t)$. Finally, we need a formula for $\widehat{\Phi_{s}}(t)$. We first note the following simple lemma, whose proof we leave to the reader.
Lemma 5.5.1. Suppose that $\Phi$ is the unramified vector in $\operatorname{Ind}_{B}^{S L_{2}} \chi_{s}$ (unnormalized induction), where

$$
\chi_{s}\left(\alpha^{\vee}(t)\right)=|t|^{1+s}
$$

and $\alpha$ is the positive root of $S L_{2}$. Then

$$
\int_{F_{v}} \Phi\left(x_{-\alpha}(r)\right) \cdot \psi(t r) d r=\frac{\zeta(s)}{\zeta(s+1) \cdot \zeta((\operatorname{ord}(t)+1) s)} .
$$

With this lemma, we can calculate $\widehat{\Phi_{s}}(t)$. We shall do so assuming that $E=F^{3}$ is split. The other cases are similarly handled. Let

$$
t=\prod_{i=0}^{3} \mu_{i}^{\prime}\left(t_{i}\right)
$$

Then

$$
\begin{aligned}
& \widehat{\Phi_{s}}(t) \\
= & \int_{F_{v}^{3}} \Phi_{s}\left(x_{-\alpha_{1}}\left(r_{1}\right) x_{-\alpha_{2}}\left(r_{2}\right) x_{-\alpha_{3}}\left(r_{3}\right) t\right) \cdot \psi\left(r_{1}+r_{2}+r_{3}\right) d r_{1} d r_{2} d r_{3} \\
= & \int_{F_{v}^{3}} \Phi_{s}\left(t \cdot\left(\prod_{i=1}^{3} x_{-\alpha_{i}}\left(t_{i}^{-1} r_{i}\right)\right) \cdot \psi\left(r_{1}+r_{2}+r_{3}\right) d r_{1} d r_{2} d r_{3}\right. \\
= & \left|t_{1}\right|^{s_{1}+1}\left|t_{2}\right|^{s_{2}+1}\left|t_{3}\right|^{s_{3}+1}\left|t_{0}\right|^{s_{1}+s_{2}+s_{3}+3} \cdot \int_{F_{v}^{3}} \Phi_{s}\left(\prod_{i=1}^{3} x_{-\alpha_{i}}\left(t_{i}^{-1} r_{i}\right)\right) \cdot \psi\left(r_{1}+r_{2}+r_{3}\right) d r_{1} d r_{2} d r_{3} \\
= & \left|t_{1}\right|^{s_{1}+2}\left|t_{2}\right|^{s_{2}+2}\left|t_{3}\right|^{s_{3}+2}\left|t_{0}\right|^{s_{1}+s_{2}+s_{3}+3} . \\
& \cdot \int_{F_{v}^{3}} \Phi_{s}\left(\prod_{i=1}^{3} x_{-\alpha_{i}}\left(r_{i}\right)\right) \cdot \psi\left(t_{1} r_{1}\right) \cdot \psi\left(t_{2} r_{2}\right) \cdot \psi\left(t_{3} r_{3}\right) d r_{1} d r_{2} d r_{3}
\end{aligned}
$$

Now, observing that

$$
\Phi_{s}\left(\alpha_{1}^{\vee}(t)\right)=|t|^{1+s_{2}+s_{3}-s_{1}} \quad \text { and so on }
$$

and appealing to Lemma 5.5.1, we obtain the following proposition.
Proposition 5.5.2. (i) When $E=F^{3}$,

$$
\begin{gathered}
\widehat{\Phi_{s}}(t)= \\
\left|t_{1}\right|^{s_{1}+2}\left|t_{2}\right|^{s_{2}+2}\left|t_{3}\right|^{s_{3}+2}\left|t_{0}\right|^{s_{1}+s_{2}+s_{3}+3} \\
\zeta\left(s_{2}+s_{3}-s_{1}\right) \\
\frac{\zeta\left(s_{2}+s_{3}-s_{1}+1\right) \cdot \zeta\left(\left(\operatorname{ord}\left(t_{1}\right)+1\right)\left(s_{2}+s_{3}-s_{1}\right)\right)}{\zeta\left(s_{3}+s_{1}-s_{2}\right)} \\
\frac{\zeta\left(s_{3}+s_{1}-s_{2}+1\right) \cdot \zeta\left(\left(\operatorname{ord}\left(t_{2}\right)+1\right)\left(s_{3}+s_{1}-s_{2}\right)\right)}{\zeta\left(s_{1}+s_{2}-s_{3}\right)} \\
\frac{\zeta\left(s_{1}+s_{2}-s_{3}+1\right) \cdot \zeta\left(\left(\operatorname{ord}\left(t_{3}\right)+1\right)\left(s_{1}+s_{2}-s_{3}\right)\right)}{2}
\end{gathered}
$$

(ii) When $E=F \times K$,

$$
\widehat{\Phi_{s}}(t)=
$$

$$
\begin{gathered}
\left|t_{1}\right|^{s_{1}+2}\left|t_{23}\right|^{2 s_{23}+4}\left|t_{0}\right|^{s_{1}+2 s_{23}+3} . \\
\frac{\zeta\left(2 s_{23}-s_{1}\right)}{\zeta\left(2 s_{23}-s_{1}+1\right) \cdot \zeta\left(\left(\operatorname{ord}\left(t_{1}\right)+1\right)\left(2 s_{23}-s_{1}\right)\right)} \\
\frac{\zeta_{K}\left(s_{1}\right)}{\zeta_{K}\left(s_{1}+1\right) \cdot \zeta_{K}\left(\left(\operatorname{ord}\left(t_{23}\right)+1\right) s_{1}\right)} .
\end{gathered}
$$

(iii) When $E$ is a field,

$$
\widehat{\Phi_{s}}(t)=\left|t_{123}\right|^{3 s+6}\left|t_{0}\right|^{3 s+3} \cdot \frac{\zeta_{E}(s)}{\zeta_{E}(s+1) \cdot \zeta_{E}\left(\left(\operatorname{ord}\left(t_{123}\right)+1\right) s\right)} .
$$

We are now ready to begin the unramified computation, Here is the main local theorem to be proved:

Theorem 5.5.3. Set

$$
Z^{*}(\varphi, \Phi, f, s)=\hat{\zeta}_{E}(\underline{s}+\underline{1}) \cdot \zeta(|\underline{s}|) \cdot Z(\varphi, \Phi, f, s) .
$$

Then

$$
Z^{*}(\varphi, \Phi, f, s)=\frac{L\left(\underline{s}, \pi, V_{E, \mathrm{Spin}}\right)}{\zeta_{E}(2 \underline{s})}
$$

Before diving into the computation, we would like to point out the three main steps in the computation. Hopefully this will make the structure of the computations more transparent.
(i) (Separation of Variables). Both sides of the identities are functions in the vector $\underline{s}$, but on the RHS, one clearly has a separation of variables. Thus, our first step is to prove that the LHS also has separation of variables. This is the most complicated step and uses the results of [BKW] regarding the decomposition of tensor products of representations of classical groups. Obviously, this step is not necessary if $E$ is a field.
(ii) (Replacing $s_{\pi}$ by $\overline{s_{\pi}}$ ). Observe that the LHS is expressed in terms of $\overline{s_{\pi}}$ via the Casselman-Shalika formula, but the RHS is in terms of $s_{\pi}$. Thus, the second step is to interpret the RHS in terms of $\overline{s_{\pi}}$. Obviously, this step is not necessary if $E=F^{3}$.
(iii) (Comparison) The final step is the comparision of the two sides. This also requires knowledge of decomposition of the tensor product of two representations. However, one of the representations will be fairly small, and so one can appeal to a more direct technique, such as Brauer's method, as opposed to using [BKW].

## 6. Unramified Computation: $E$ a field.

As might be expected, once the subtleties of non-split groups are understood, the computation there turns out to be simpler than the split case. Thus we begin with the case when $E$ is a field.
6.1. Zeta integral side. Recall that we have the two versions of Satake parameter:

$$
s_{\pi} \in H_{E}^{\prime} \cong G_{2}(\mathbb{C}) \quad \text { and } \quad \overline{s_{\pi}} \in H_{E}^{\vee} \cong G_{2}(\mathbb{C}) .
$$

If we set

$$
x=q^{-3 s},
$$

then the LHS of Theorem 5.5.3 is equal to:

$$
\begin{aligned}
& (1-x)^{-1} \cdot \sum_{k_{0}, k} \epsilon_{3}\left(k_{0}+1\right) \cdot \operatorname{Tr}\left(\overline{s_{\pi}} \mid V_{k_{0}, k}^{H_{E}^{\vee}}\right) \cdot x^{k_{0}+k} \cdot \frac{1-x^{k+1}}{1-x} \\
& =(1-x)^{-1} \cdot \sum_{k_{0}, k} \sum_{l=0}^{k} \epsilon_{3}\left(k_{0}+1\right) \cdot \operatorname{Tr}\left(\overline{s_{\pi}} \mid V_{k_{0}, k}^{H_{\bigvee}^{\vee}}\right) \cdot x^{l+k_{0}+k} .
\end{aligned}
$$

Here, $V_{k_{0}, k}$ is the irreducible representation of $H_{E}^{\vee}$ with highest weight $\left(k_{0}, k\right)$, where $(1,0)$ stands for the fundamental weight attached to the 7-dim representation of $H_{E}^{\vee}$ and $(0,1)$ that of the adjoint representation.
6.2. $L$-function side. On the other hand, to explicate the other side of the main theorem, which involves the Spin L-function, we consider the restriction of $V_{E, \text { Spin }}$ to the subgroup $H_{E}^{\prime} \times \operatorname{Gal}(\bar{F} / F)$. This decomposes as:

$$
V_{E, \text { Spin }} \cong\left(V_{1,0}^{H_{E}^{\prime}} \oplus \mathbf{1}\right) \boxtimes\left(\mathbf{1} \oplus \chi_{E} \oplus \chi_{E}^{-1}\right)
$$

where $V_{1,0}^{H_{E}^{\prime}}$ is the 7-dimensional fundamental representation of $H_{E}^{\prime} \cong G_{2}(\mathbb{C})$. Since $t_{\pi}=$ ( $s_{\pi}$, Frob) with $s_{\pi} \in H_{E}^{\prime}$, we see that

$$
\operatorname{det}\left(1-q^{-s} t_{\pi} \mid V_{E, \text { Spin }}\right)=\left(1-q^{-3 s}\right) \cdot \operatorname{det}\left(1-q^{-3 s} s_{\pi} \mid V_{1,0}^{H_{E}^{\prime}}\right)
$$

Thus the RHS of the identity in Theorem 5.5.3 is equal to

$$
\begin{aligned}
& \frac{L\left(s, \pi, V_{E, \text { Spin }}\right)}{\zeta_{E}(2 s) \cdot \zeta(3 s)} \\
= & \left(1-q^{-6 s}\right) \cdot\left(1-q^{-3 s}\right)^{-1} \cdot \operatorname{det}\left(1-q^{-3 s} s_{\pi}^{3} \mid V_{1,0}^{H_{E}^{\prime}}\right)^{-1} \\
= & (1-x)^{-1} \cdot\left(1-x^{2}\right) \cdot \sum_{n} x^{n} \cdot \operatorname{Tr}\left(s_{\pi}^{3} \mid S_{y m}{ }^{n}\left(V_{1,0}^{H_{E}^{\prime}}\right)\right) \\
= & (1-x)^{-1} \cdot \sum_{n} \operatorname{Tr}\left(s_{\pi}^{3} \mid V_{n, 0}^{H_{E}^{\prime}}\right) \cdot x^{n} .
\end{aligned}
$$

In the last equality above, we have used the fact that

$$
\operatorname{Sym}^{n}\left(V_{1,0}\right) \cong \operatorname{Sym}^{n-2}\left(V_{1,0}\right) \oplus V_{n, 0} .
$$

For this fact, see the table on [ $\mathrm{Br}, \mathrm{Pg}$. 13].
6.3. Comparison. We can now attempt to compare the two sides. Unfortunately, one side of the identity is expressed in terms of $\overline{s_{\pi}}$ while the other in terms of $s_{\pi}$. Thus, we need to express the trace of $s_{\pi}^{3}$ on $V_{n, 0}^{H_{E}^{\prime}}$ in terms of the trace of $\overline{s_{\pi}}$ on some representation of $H_{E}^{\vee}$. Using the explicit description of the map $\iota^{*}$ in (5.3) and the Weyl character formula, we obtain:

## Lemma 6.3.1.

$$
\operatorname{Tr}\left(s_{\pi} \mid V_{n, 0}^{H_{E}^{\prime}}\right)=\frac{\operatorname{Tr}\left(\overline{s_{\pi}} \mid V_{2, n}^{H_{V}^{\vee}}\right)}{\operatorname{Tr}\left(\overline{s_{\pi}} \mid V_{2,0}^{H_{E}^{\vee}}\right)} .
$$

If we cancel the factor $(1-x)^{-1}$ from both sides of the desired identity and compare the coefficient of $x^{n}$, we see that to complete the proof of Theorem 5.5.3 in this case, we need to prove:

## Proposition 6.3.2.

$$
\frac{\operatorname{Tr}\left(\overline{s_{\pi}} \mid V_{2, n}\right)}{\operatorname{Tr}\left(\overline{s_{\pi}} \mid V_{2,0}\right)}=\sum_{k+k_{0} \leq n \leq 2 k+k_{0}} \epsilon_{3}\left(k_{0}+1\right) \cdot \operatorname{Tr}\left(\overline{s_{\pi}} \mid V_{k_{0}, k}\right) .
$$

The rest of this section is devoted to the proof of this proposition.
6.4. Reduction to $S L_{3}$. The long root subgroups of $G_{2}(\mathbb{C})$ generate a subgroup of $G_{2}$ isomorphic to $S L_{3}(\mathbb{C})$, sharing the same maximal torus $S_{E}^{\vee}$. The roots $\{\beta, 3 \alpha+\beta\}$ of $G_{2}$ form a system of simple roots for $S L_{3}$, whose fundamental weights are

$$
\left\{\begin{array}{l}
\omega_{1}=2 \alpha+\beta \\
\omega_{2}=\alpha+\beta
\end{array}\right.
$$

We write $W_{a, b}$ for the irreducible representation of $S L_{3}(\mathbb{C})$ with highest weight $a \omega_{1}+b \omega_{2}$ and let $\chi_{a, b}$ be its character. Henceforth, we shall use the coordinates given by the weights $\omega_{1}$ and $\omega_{2}$ of $S L_{3}$.

Using the Weyl character formula, one checks that

$$
\operatorname{Tr}\left(s \mid V_{a, b}\right)=\frac{\operatorname{Tr}\left(s \mid W_{a+b+1, b}\right)-\operatorname{Tr}\left(s \mid W_{b, a+b+1}\right)}{\operatorname{Tr}\left(s \mid W_{1,0}\right)-\operatorname{Tr}\left(s \mid W_{0,1}\right)} .
$$

Thus the identity of Prop. 6.3.2 is:

$$
\begin{gather*}
\left(\chi_{n+3, n}-\chi_{n, n+3}\right)\left(\chi_{1,0}-\chi_{0,1}\right)  \tag{6.4.1}\\
=\left(\sum_{k+k_{0} \leq n \leq 2 k+k_{0}} \epsilon_{3}\left(k_{0}+1\right) \cdot\left(\chi_{k+k_{0}+1, k}-\chi_{k, k+k_{0}+1}\right)\right)\left(\chi_{3,0}-\chi_{0,3}\right) .
\end{gather*}
$$

We use Brauer's method, which for these purposes is conveniently stated as follows. Let $\xi$ be any virtual character of $S L_{3}(\mathbb{C})$, given by

$$
\xi(t)=\sum_{\nu} m(\nu) t^{\nu},
$$

with the sum being over the weight lattice, $m(\nu) \in \mathbb{Z}$ for each $\nu$, and $m(\nu)=0$ for almost all $\nu$. Note also that, since $\chi$ is a virtual character, $m(w \nu)=m(\nu)$ for any weight $\nu$ and any $w$ in the Weyl group. For $\lambda$ any weight, let

$$
A_{\lambda}(t)=\sum_{w \in W_{S L_{3}}} \operatorname{sgn}(w) t^{w \lambda}
$$

and let $\chi_{\lambda}=A_{\lambda+\rho} / A_{\rho}$. Then

$$
\begin{equation*}
\xi \chi_{\lambda}=\sum_{\nu} m(\nu) \chi_{\lambda+\nu} \tag{6.4.2}
\end{equation*}
$$

Observe that $\chi_{\lambda}$, so defined, is equal to the character of the irreducible representation with highest weight $\lambda$ when $\lambda$ is dominant, and is zero when $\lambda+\rho$ has a stabilizer in the Weyl group, and that $\chi_{\lambda}+\chi_{\mu}=0$ whenever $\lambda+\rho$ and $\mu+\rho$ are related by a simple reflection.

We apply (6.4.2) to both sides of (6.4.1), and summarize the answer by a picture of the weight lattice where each node $\tau$ is labeled with the multiplicity of the corresponding $\chi_{\tau}$. As it turns out all multiplicities are 0,1 or -1 , and 0 's are omitted. The application to the left hand side, with $\xi=\chi_{1,0}-\chi_{0,1}$, is straightforward, and for $n \geq 1$ yields two hexagons of side 1 , centered at $(n+3, n)$ and $(n, n+3)$ (here, and throughout, weights are expressed in terms of the basis of fundamental weights) with alternating signs. See Figure 1.


Figure 1. the LHS of (6.4.1).
Next we apply (6.4.2) to the right hand side of (6.4.1) with $\xi=\chi_{3,0}-\chi_{0,3}$. The representations $W_{3,0}$ and $W_{0,3}$ have nine weights each, but six of them are common to both, and cancel, leaving a hexagon with sides of length three, and signs alternating. As may be expected, the details are a bit different for the first few values of $n$. We give an argument which works for $n \geq 6$. The reader may find it an enjoyable diversion to work out the remaining cases. Alternatively, one could use the computer package $\operatorname{LiE}$ to verify the cases $n<6$.

Now, we note that

$$
\begin{aligned}
& \sum_{k+k_{0} \leq n \leq 2 k+k_{0}} \epsilon_{3}\left(k_{0}+1\right) \cdot\left(\chi_{k+k_{0}+1, k}-\chi_{k, k+k_{0}+1}\right) \\
= & \sum_{0 \leq b<a \leq n+1 \leq a+b} \epsilon_{3}(a-b) \chi_{a, b}-\sum_{0 \leq a<b \leq n+1 \leq a+b} \epsilon_{3}(b-a) \chi_{a, b} \\
= & \sum_{0 \leq a, b \leq n+1 \leq a+b} \epsilon_{3}(a-b) \chi_{a, b} .
\end{aligned}
$$



Figure 2
The sum is now over the lattice points contained in the triangle with vertices $(n+1,0),(0, n+1)$ and ( $n+1, n+1$ ). Let us denote this triangle by $\Delta_{n+1}$.

We now multiply the above sum by $\xi$, using (6.4.2) on each term. It's clear that $\chi_{\tau}$ appears in the product with nonzero multiplicity only if $\tau$ lies in one of the six translates of $\Delta_{n+1}$ by the weights $\nu$ such that $m(\nu) \neq 0$. Setting

$$
\epsilon_{3}(\tau)=\epsilon_{3}(a-b) \quad \text { for } \tau=(a, b)
$$

observe that

$$
\epsilon_{3}(\tau)=\epsilon_{3}(\tau+\nu) \quad \text { for each such } \nu
$$

Figure 2 shows $\Delta_{n+1}$, along with its six translates, visualized as pointing straight up. With things drawn this way, $\chi(\tau)$ is constant on vertical lines, and $\rho$ points straight up. The multiplicity with which $\chi_{\tau}=\chi_{a, b}$ appears in the product is

$$
\epsilon_{3}(\tau) \sum_{\nu: \tau-\nu \in \Delta_{n+1}} m(\nu) .
$$

Now, we must analyze the cancellation coming from the overlap of the six translates of this triangle.

This is displayed in Figure 3: we first see how many of the translates with $m(\nu)=1$ each point lies in, and then how many of the ones with $m(\nu)=-1$. Subtracting, we see that the only weights which contribute lie in six rhombi with side 2 , one on each side of each vertex of our triangle. (The precise picture is only valid for $n$ at least 6 but the reader will find that this description holds for $n$ as small as 2.)

We consider the two rhombi near the vertex $(n+1, n+1)$. As we see in Figure 4, the points on the long diagonals of our two rhombi do not contribute because $\epsilon_{3}(a-b)$ is zero on these lines. The remaining points give precisely the two hexagons of side 1 we saw before, and the signs match.

Next we turn to the four rhombi at the vertices $(n+1,0)$ and $(0, n+1)$. In Figure 5 one of these rhombi is displayed in grey, while its translate by $\rho$ is shown in black. We see that,


Figure 3. Positive and negative translates of $\Delta_{n+1}$ cancel almost completely.


Figure 4. grey numbers indicate $\epsilon_{3}$; black signs indicate $m(\nu)$; cf Fig. 1.
upon translation by $\rho$, the three points on the short diagonal are taken to points which are stabilized by an elementary reflection, so their contributions vanish. The remaining six points are related by this elementary reflection, so their contributions cancel. This last statement relies on the fact that the values of $\epsilon_{3}(a-b)$ at corresponding points are equal, which may also be observed from Figure 5.

A similar consideration shows the desired vanishing for the other 3 rhombi. Proposition 6.3.2 is proved.
7. Unramified Computation: $E=F \times K$

Now we come to the case when $E=F \times K$. In this case, the element $t$ of $S_{E}(F)$ is of the form

$$
t=\mu_{0}^{\prime}\left(t_{0}\right) \mu_{1}^{\prime}\left(t_{1}\right)\left(\mu_{2}^{\prime} \mu_{3}^{\prime}\right)\left(t_{23}\right)=\mu_{25}^{\prime}\left(\varpi^{k_{0}}\right) \mu_{1}^{\prime}\left(\varpi^{k_{1}}\right)\left(\mu_{2}^{\prime} \mu_{3}^{\prime}\right)\left(\varpi^{k_{23}}\right)
$$



Figure 5
We recall the Casselman-Shalika formula:

$$
\delta_{B_{E}^{\prime}}(t)^{-1 / 2} W_{V_{E}^{\prime}, \psi}(t \cdot \varphi)=\operatorname{Tr}\left(\overline{s_{\pi}} \mid V_{k_{1}, k_{0}, k_{23}}^{H_{E}^{\vee}}\right)
$$

Here, $V_{k_{1}, k_{0}, k_{23}}^{H_{E}^{\vee}}$ denotes the irreducible representation of $H_{E}^{\vee} \cong S p_{6}(\mathbb{C})$ with highest weight $k_{1} \nu_{1}+k_{0} \nu_{2}+k_{23} \nu_{3}$ where $\nu_{1}, \nu_{2}, \nu_{3}$ are the fundamental weights of $S p_{6}(\mathbb{C})$ numbered such that $\nu_{1}$ corresponds to the standard representation.
7.1. Zeta integral side. Let $x=q^{-s_{1}}$ and $y=q^{-2 s_{23}}$. Then the LHS of Theorem 5.5.3 is:
$(1-x y)^{-1} \cdot \sum_{k_{i}=0}^{\infty} x^{k_{0}+k_{1}} y^{k_{0}+k_{23}} \epsilon_{2}\left(k_{0}+1\right) \cdot(-1)^{k_{1}} \cdot \frac{1-\left(x^{-1} y\right)^{k_{1}+1}}{1-x^{-1} y} \frac{1-x^{2\left(k_{23}+1\right)}}{1-x^{2}} \operatorname{Tr}\left(\overline{s_{\pi}} \mid V_{k_{1}, k_{0}, k_{23}}^{S p_{6}}\right)$.
Next we have the following crucial lemma, which gives a separation of variables in the above sum:

Lemma 7.1.1. The LHS of Theorem 5.5.3 is equal to:

$$
\left(\sum_{k=0}^{\infty}(-x)^{k} \operatorname{Tr}\left(\overline{s_{\pi}} \mid V_{k, 0,0}^{S p_{6}}\right)\right)\left(\sum_{\ell, m=0}^{\infty} y^{\ell+m}(-1)^{\ell} \operatorname{Tr}\left(\overline{s_{\pi}} \mid V_{\ell, 0, m}^{S p_{6}}\right)\right)
$$

So as not to disrupt the flow of the argument, we defer the proof of the lemma to the end of the section.
7.2. $L$-function side. Now we consider the $L$-function side of the main theorem, which is a priori expressed in terms of $t_{\pi} \in{ }^{L} G_{E}$. We know that $V_{E, \text { Spin }}$ is reducible as a representation of ${ }^{L} G_{E}$ : it decomposes as:

$$
V_{E, \text { Spin }}=V_{\mu_{1}} \bigoplus\left(V_{\mu_{2}} \oplus V_{\mu_{3}}\right)
$$

As in the previous section, we consider the restriction of $V_{E, \text { Spin }}$ to the subgroup

$$
H_{E}^{\prime} \times \operatorname{Gal}(\bar{F} / F)=\operatorname{Spin}_{7}(\mathbb{C}) \times \operatorname{Gal}(\bar{F} / F) .
$$

It is not difficult to see that:

$$
V_{\mu_{1}} \cong\left(V_{1,0,0}^{H_{E}^{\prime}} \boxtimes \chi_{K}\right) \bigoplus(\mathbf{1} \boxtimes \mathbf{1})
$$

where $V_{1,0,0}^{H_{E}^{\prime}}$ is the 7 -dimensional standard representation of $S p i i_{7}$, and

$$
V_{\mu_{2}} \oplus V_{\mu_{3}} \cong V_{0,0,1}^{H_{E}^{\prime}} \boxtimes\left(\mathbf{1} \oplus \chi_{K}\right)
$$

where $V_{0,0,1}^{H_{E}^{\prime}}$ is the 8-dimensional Spin representation of $\operatorname{Spin}_{7}$. Thus, we obtain:

Next we need to express the above in terms of $\overline{s_{\pi}}$ rather than $s_{\pi}$. For this, we need to know explicitly the map

$$
\iota^{*}: T_{E}^{\vee} \longrightarrow S_{E}^{\vee}
$$

The maximal torus $T_{E}^{\vee}$ of $\operatorname{Spin}_{8}(\mathbb{C})$ is conveniently parametrized by five complex variables $\left(u, t_{1}, t_{2}, t_{3}, t_{4}\right)$ subject to the relation $t_{1} t_{2} t_{3} t_{4}=u^{2}$. With this parametrization, the map to the maximal torus $S_{E}^{\vee}$ of $S p_{6}(\mathbb{C})$ is simply

$$
\left(u, t_{1}, t_{2}, t_{3}, t_{4}\right) \mapsto\left(t_{1}, t_{2}, t_{3}\right) .
$$

From this, it is not hard to see that

$$
\operatorname{det}\left(1+q^{-s_{1}} s_{\pi} \mid V_{1,0,0}^{S p_{i n}}\right)=\left(1+q^{-s_{1}}\right) \cdot \operatorname{det}\left(1+q^{-s_{1}} \bar{s}_{\pi} \mid V_{1,0,0}^{S p_{6}}\right),
$$

where $V_{1,0,0}^{S p_{6}}$ is the 6 -dimensional standard representation of $S p_{6}(\mathbb{C})$. Thus, we have:

$$
\begin{aligned}
& \frac{1}{\left(1+q^{-s_{1}}\right)} \cdot \frac{1}{\operatorname{det}\left(1+q^{\left.-s_{1} s_{\pi} \mid V_{1,0,0}^{S p i n 7}\right)}\right.} \\
= & \zeta\left(2 s_{1}\right) \cdot \sum_{k=0}^{\infty} \operatorname{Tr}\left(\bar{s}_{\pi} \mid S y m^{k} V_{1,0,0}^{S p_{6}}\right) \cdot(-1)^{k} \cdot q^{-k s_{1}} \\
= & \zeta\left(2 s_{1}\right) \cdot \sum_{k=0}^{\infty} \operatorname{Tr}\left(\bar{s}_{\pi} \mid V_{k, 0,0}^{S p_{6}}\right) \cdot(-1)^{k} \cdot q^{-k s_{1}},
\end{aligned}
$$

where for the last equality, we have used the well-known fact that

$$
\operatorname{Sym}^{n}\left(V_{1,0,0}^{S p_{6}}\right)=V_{n, 0,0}^{S p_{6}}
$$

which can be found in the table on [ $\mathrm{Br}, \mathrm{Pg}$. 13].
On the other hand, by [Br, Pg. 13] again, we see that

$$
\operatorname{det}\left(1-q^{-2 s_{23}} s_{\pi}^{2} \mid V_{0,0,1}^{S_{p i i_{7}}}\right)^{-1}=\zeta_{K}\left(2 s_{23}\right) \sum_{n=0}^{\infty} q^{-2 n s_{23}} \operatorname{Tr}\left(s_{\pi}^{2} \mid V_{0,0, n}^{\text {Spin }_{7}}\right),
$$

where $V_{0,0, n}^{\text {Spin }_{7}}$ denotes the irreducible representation of $\operatorname{Spin}_{7}$ with highest weight equal to $n$ times the third fundamental weight. We claim:

## Lemma 7.2.1.

$$
\operatorname{Tr}\left(s_{\pi}^{2} \mid V_{0,0, n}^{S_{p i n}}\right)=\frac{\operatorname{Tr}\left(\bar{s}_{\pi} \mid V_{1,1, n}^{S p_{6}}\right)}{\operatorname{Tr}\left(\bar{s}_{\pi} \mid V_{1,1,0}^{S p_{6}}\right)} .
$$

Proof. Let us view $s_{\pi}=\left(u, t_{1}, t_{2}, t_{3}, t_{4}\right)$ as above. The fact that $s_{\pi}$ actually lies in the maximal torus of $H_{E}^{\prime} \cong \operatorname{Spin}_{7}(\mathbb{C})$ implies that $t_{4}=1$. The coordinates of $s_{\pi}^{2}$ are then
$\left(t_{1} t_{2} t_{3}, t_{1}^{2}, t_{2}^{2}, t_{3}^{2}, 1\right)$. Let $s(t)=t-t^{-1}$. Then the Weyl character formula can be given quite nicely in terms of determinants in this case:

$$
\operatorname{Tr}\left(s_{\pi}^{2} \mid V_{0,0, n}^{S p i n_{7}}\right)=\frac{\left|\begin{array}{lll}
s\left(t_{1}^{n+5}\right) & s\left(t_{2}^{n+5}\right) & s\left(t_{3}^{n+5}\right) \\
s\left(t_{1}^{n+3}\right) & s\left(t_{2}^{n+3}\right) & s\left(t_{3}^{n+3}\right) \\
s\left(t_{1}^{n+1}\right) & s\left(t_{2}^{n+1}\right) & s\left(t_{3}^{n+1}\right)
\end{array}\right|}{\left|\begin{array}{lll}
s\left(t_{1}^{5}\right) & s\left(t_{2}^{5}\right) & s\left(t_{3}^{5}\right) \\
s\left(t_{1}^{3}\right) & s\left(t_{2}^{3}\right) & s\left(t_{3}^{3}\right) \\
s\left(t_{1}^{1}\right) & s\left(t_{2}^{1}\right) & s\left(t_{3}^{1}\right)
\end{array}\right|}
$$

A similar formula holds for the characters of $S p_{6}(\mathbb{C})$ :

$$
\operatorname{Tr}\left(\bar{s}_{\pi} \mid V_{k_{1}, k_{2}, k_{3}}^{S p_{6}}\right)=\frac{\left|\begin{array}{ccc}
s\left(t_{1}^{k_{1}+k_{2}+k_{3}+3}\right) \\
s\left(t_{1}^{k_{2}+k_{3}+2}\right) & s\left(t_{2}^{k_{1}+k_{2}+k_{3}+3}\right) & s\left(t_{3}^{k_{1}+k_{2}+k_{3}+3}\right) \\
s\left(t_{1}^{k_{3}+1}\right) & s\left(t_{2}^{k_{2}+k_{3}+2}\right) & s\left(t_{3}^{k_{2}+k_{3}+2}\right) \\
s\left(t_{2}^{k_{3}+1}\right) & s\left(t_{3}^{k_{3}+1}\right)
\end{array}\right|}{}\left|\begin{array}{ccc}
s\left(t_{1}^{3}\right) & s\left(t_{2}^{3}\right) & s\left(t_{3}^{3}\right) \\
s\left(t_{2}^{2}\right) & s\left(t_{2}^{2}\right) & s\left(t_{3}^{2}\right) \\
s\left(t_{1}^{1}\right) & s\left(t_{2}^{1}\right) & s\left(t_{3}^{1}\right)
\end{array}\right|
$$

So we see that

$$
\operatorname{Tr}\left(s_{\pi}^{2} \mid V_{0,0, n}^{S p i n}\right)=\frac{\operatorname{Tr}\left(\bar{s}_{\pi} \mid V_{1,1, n}^{S p_{6}}\right)}{\operatorname{Tr}\left(\bar{s}_{\pi} \mid V_{1,1,0}^{S p_{6}}\right)}
$$

as desired.
7.3. Comparison. Summarizing, we have shown that:

$$
\frac{L\left(s_{1}, \pi, V_{\mu_{1}}\right)}{\zeta\left(2 s_{1}\right)}=\left(\sum_{k=0}^{\infty}(-x)^{k} \operatorname{Tr}\left(\bar{s}_{\pi} \mid V_{k, 0,0}^{S p_{6}}\right)\right)
$$

and

$$
\frac{L\left(s_{23}, \pi, V_{\mu_{2}} \oplus V_{\mu_{3}}\right)}{\zeta_{K}\left(2 s_{23}\right)}=\left(\sum_{n=0}^{\infty} y^{n} \cdot \frac{\operatorname{Tr}\left(\bar{s}_{\pi} \mid V_{1,1, n}^{S p_{6}}\right)}{\operatorname{Tr}\left(\bar{s}_{\pi} \mid V_{1,1,0}^{S p_{6}}\right)}\right) .
$$

Thus the RHS of the identity in Theorem 5.5.3 is:

$$
\left(\sum_{k=0}^{\infty}(-x)^{k} \operatorname{Tr}\left(\bar{s}_{\pi} \mid V_{k, 0,0}^{S p_{6}}\right)\right) \cdot\left(\sum_{n=0}^{\infty} y^{n} \cdot \frac{\operatorname{Tr}\left(\bar{s}_{\pi} \mid V_{1,1, n}^{S p_{6}}\right)}{\operatorname{Tr}\left(\bar{s}_{\pi} \mid V_{1,1,0}^{S p_{6}}\right)}\right) .
$$

Comparing this with Lemma 7.1.1, we see that it remains to prove:

## Proposition 7.3.1.

$$
\frac{\operatorname{Tr}\left(s \mid V_{1,1, n}^{S p_{6}}\right)}{\operatorname{Tr}\left(s \mid V_{1,1,0}^{S S_{6}}\right)}=\sum_{\ell=0}^{n}(-1)^{\ell} \cdot \operatorname{Tr}\left(s \mid V_{\ell, 0, n-\ell}^{S p_{6}}\right) .
$$

Proof. This can be proved using Brauer's method, as in Prop. 6.3.2, but is harder to explain in this case since we will be working in 3-dimensional space. We give a proof using the expression of the characters in determinantal form.

The identity to be proven can be restated as

$$
\begin{aligned}
& =\left|\begin{array}{lll}
s\left(t_{1}^{n+5}\right) & s\left(t_{2}^{n+5}\right) & s\left(t_{3}^{n+5}\right) \\
s\left(t_{1}^{n+3}\right) & s\left(t_{2}^{n+3}\right) & s\left(t_{3}^{n+3}\right) \\
s\left(t_{1}^{n+1}\right) & s\left(t_{2}^{n+1}\right) & s\left(t_{3}^{n+1}\right)
\end{array}\right| .
\end{aligned}
$$

We first check that

$$
\Delta(t):=\frac{\left|\begin{array}{lll}
s\left(t_{1}^{5}\right) & s\left(t_{2}^{5}\right) & s\left(t_{3}^{5}\right) \\
s\left(t_{3}^{3}\right) & s\left(t_{2}^{3}\right) & s\left(t_{3}^{3}\right) \\
s\left(t_{1}^{1}\right) & s\left(t_{2}^{1}\right) & s\left(t_{3}^{1}\right)
\end{array}\right|}{\left|\begin{array}{lll}
s\left(t_{1}^{3}\right) & s\left(t_{2}^{3}\right) & s\left(t_{3}^{3}\right) \\
s\left(t_{1}^{2}\right) & s\left(t_{2}^{2}\right) & s\left(t_{3}^{2}\right) \\
s\left(t_{1}^{1}\right) & s\left(t_{2}^{1}\right) & s\left(t_{3}^{1}\right)
\end{array}\right|}=\left(t_{1}^{-1}+t_{2}^{-1}\right)\left(t_{1}^{-1}+t_{3}^{-1}\right)\left(t_{2}^{-1}+t_{3}^{-1}\right)\left(1+t_{1} t_{2}\right)\left(1+t_{1} t_{3}\right)\left(1+t_{2} t_{3}\right)
$$

Next we note that

$$
\left|\begin{array}{ccc}
s\left(t_{1}^{n+3}\right) & s\left(t_{2}^{n+3}\right) & s\left(t_{3}^{n+3}\right)  \tag{7.3.2}\\
s\left(t_{1}^{m+2}\right) & s\left(t_{2}^{m+2}\right) & s\left(t_{3}^{m+2}\right) \\
s\left(t_{1}^{m+1}\right) & s\left(t_{2}^{m+1}\right) & s\left(t_{3}^{m+1}\right)
\end{array}\right|=\sum_{\epsilon \in\{ \pm 1\}^{3}}\left|\begin{array}{ccc}
\epsilon_{1} t_{1}^{\epsilon_{1}(n+3)} & \epsilon_{2} t_{2}^{\epsilon_{2}(n+3)} & \epsilon_{3} t_{3}^{\epsilon_{3}(n+3)} t_{1}^{\epsilon_{1}(m+2)} \\
\epsilon_{2} t_{2}^{\epsilon_{2}(m+2)} & \epsilon_{3} t_{3}(m+2) \\
\epsilon_{1} t_{1}^{\epsilon_{1}(m+1)} & \epsilon_{2} t_{2}^{\epsilon_{2}(m+1)} & \epsilon_{3} t_{3}^{\epsilon_{3}(m+1)}
\end{array}\right| .
$$

Now expand each of these determinants by minors on the first row, we see that this is equal to

$$
\begin{aligned}
& \sum_{\epsilon \in\{ \pm 1\}^{3}} \sum_{i \in \mathbb{Z} / 3 \mathbb{Z}} \epsilon_{i} \epsilon_{i}^{\epsilon_{i} \cdot(n+3)} \cdot\left|\begin{array}{ll}
\epsilon_{i+1} t_{i+1}^{\epsilon_{i+1}(m+2)} & \epsilon_{i+2} t_{i+2}^{\epsilon_{i+2}(m+2)} \\
\epsilon_{i+1} t_{i+1}^{\epsilon_{i+1}}(m+1) & \epsilon_{i+2} t_{i+2}(m+1)
\end{array}\right| \\
= & \sum_{\epsilon \in\{ \pm 1\}^{3}} \sum_{i \in \mathbb{Z} / 3 \mathbb{Z}} \epsilon_{i} \epsilon_{i}^{\epsilon_{i} \cdot(n+3)} \cdot\left|\begin{array}{ll}
\epsilon_{i+1} t_{i+1}^{2 \epsilon_{i+1}} & \epsilon_{i+2} t_{i+1}^{2 \epsilon_{i+2}} \\
\epsilon_{i+1} t_{i+1}^{\epsilon_{i+1}} & \epsilon_{i+2} t_{i+2}^{i+2}
\end{array}\right| \cdot t_{i+1}^{\epsilon_{i+1} m} t_{i+2}^{\epsilon_{i+2} m}
\end{aligned}
$$

Now summing over $m$, we see that

$$
\begin{align*}
& \sum_{m=0}^{n}(-1)^{n-m} \cdot\left|\begin{array}{ccc}
s\left(t_{1}^{n+3}\right) & s\left(t_{2}^{n+3}\right) & s\left(t_{3}^{n+3}\right) \\
s\left(t_{1}^{m+2}\right) & s\left(t_{2}^{m+2}\right) & s\left(t_{3}^{m+2}\right) \\
s\left(t_{1}^{m+1}\right) & s\left(t_{2}^{m+1}\right) & s\left(t_{3}^{m+1}\right)
\end{array}\right| \\
= & \sum_{\epsilon \in\{ \pm 1\}^{3}} \sum_{i \in \mathbb{Z} / 3 \mathbb{Z}}(-1)^{n} \cdot \epsilon_{1} \epsilon_{2} \epsilon_{3} \cdot t_{i}^{(n+3) \epsilon_{i}} t_{i+1}^{\epsilon_{i+1}} t_{i+2}^{\epsilon_{i+2}} \cdot\left(t_{i+1}^{\epsilon_{i+1}}-t_{i+2}^{\epsilon_{i+2}}\right) \cdot \frac{1-\left(-t_{i+1}^{\epsilon_{i+1}} t_{i+2}^{\epsilon_{i+2}}\right)^{n+1}}{1+t_{i+1}^{\epsilon_{i+1}} t_{i+2}^{\epsilon_{i+2}}} . \tag{7.3.3}
\end{align*}
$$

If we let $T_{i}=t_{i}^{n}$, then this last expression is a linear combination $T_{i}^{\epsilon_{i}}$ and $T_{1}^{\epsilon_{1}} T_{2}^{\epsilon_{2}} T_{3}^{\epsilon_{3}}$ with coefficients given by rational functions of $t_{i}$ 's essentially independent of $n$ (the only dependence on $n$ being the factor $\left.(-1)^{n}\right)$.

Similarly, we expand the RHS of the desired identity in a manner analogous to (7.3.2). It is equal to:

$$
\sum_{\epsilon \in\{ \pm 1\}^{3}} \epsilon_{1} \epsilon_{2} \epsilon_{3} T_{1}^{\epsilon_{1}} T_{2}^{\epsilon_{2}} T_{3}^{\epsilon_{3}} \cdot\left|\begin{array}{ccc}
t_{1}^{5 \epsilon_{1}} & t_{2}^{5 \epsilon_{2}} & t_{3}^{5 \epsilon_{3}} \\
t_{1}^{3 \epsilon_{1}} & t_{2}^{3 \epsilon_{2}} & t_{3}^{3 \epsilon_{3}} \\
t_{1}^{\epsilon_{1}} & t_{2}^{\epsilon_{2}} & t_{3}^{\epsilon_{3}}
\end{array}\right| .
$$

Treat both sides of the desired identity as polynomials in $T_{1}, T_{2}, T_{3}, T_{1}^{-1}, T_{2}^{-1}, T_{3}^{-1}$, with coefficients in $\mathbb{C}\left(t_{1}, t_{2}, t_{3}\right)$. Then we need to show:

- the coefficient of $T_{i}^{\epsilon_{i}}$ on the LHS vanishes. By symmetry, it suffices to show

$$
\sum_{\epsilon \in\{ \pm 1\}^{2}} \frac{\epsilon_{1} \epsilon_{2} t_{1}^{\epsilon_{1}} t_{2}^{\epsilon_{2}} \cdot\left(t_{1}^{\epsilon_{1}}-t_{2}^{\epsilon_{2}}\right)}{1+t_{1}^{\epsilon_{1}} t_{2}^{\epsilon_{2}}}=0
$$

This is immediately verifiable by hand.

- the coefficient of $T_{1}^{\epsilon_{1}} T_{2}^{\epsilon_{2}} T_{3}^{\epsilon_{3}}$ on both sides agree. By symmetry again, it suffices to verify:

$$
\left(\frac{t_{1}^{3} t_{2}^{2} t_{3}^{2}\left(t_{2}-t_{3}\right)}{1+t_{2} t_{3}}+\frac{t_{2}^{3} t_{3}^{2} t_{1}^{2}\left(t_{3}-t_{1}\right)}{1+t_{3} t_{1}}+\frac{t_{3}^{3} t_{1}^{2} t_{2}^{2}\left(t_{1}-t_{2}\right)}{1+t_{1} t_{2}}\right) \cdot \Delta(t)=\left|\begin{array}{ccc}
t_{1}^{5} & t_{2}^{5} & t_{3}^{5} \\
t_{1}^{3} & t_{2}^{3} & t_{3}^{3} \\
t_{1}^{1} & t_{2}^{1} & t_{3}^{1}
\end{array}\right| .
$$

Again, this is easily verified by hand.
The proposition is proved.
7.4. Proof of Lemma 7.1.1. Finally, we need to give the proof of Lemma 7.1.1.

For this, we use a formula which goes back to Murnaghan and Littlewood, and has been given a nice interpretation in terms of "universal characters" in [Ko] and [K-T]. This is discussed in greater detail in the appendix. We state it in the notation of the appendix. We associate the highest weight of $V_{k_{1}, k_{2}, k_{3}}^{S P_{6}}$ with the partition $\left(k_{1}+k_{2}+k_{3}, k_{2}+k_{3}, k_{3}\right)$, and for $\lambda$ a partition with at most 3 parts, let $\chi_{S p_{6}}(\lambda)$ denote the character of the representation with the corresponding highest weight. We consider also finite formal sums of partitions, and extend $\chi_{S p_{6}}$ to these by $\mathbb{Z}$-linearity.

Thus, the right hand side of the identity is

$$
\begin{equation*}
(1-x y)\left(\sum_{k=0}^{\infty}(-x)^{k} \chi_{S p_{6}}(k)\right)\left(\sum_{0 \leq \mu_{2} \leq \mu_{1}} y^{\mu_{1}}(-1)^{\mu_{1}-\mu_{2}} \chi_{S p_{6}}\left(\mu_{1} \mu_{2}^{2}\right)\right) . \tag{7.4.1}
\end{equation*}
$$

On the set of all partitions, we define operations • and / by

$$
\lambda \cdot \mu=\sum_{\nu} L R_{\lambda, \mu}^{\nu} \nu
$$

and

$$
\lambda / \mu=\sum_{\nu} L R_{\mu, \nu}^{\lambda} \nu,
$$

where the Littlewood-Richardson coefficient $L R_{\mu, \nu}^{\lambda}$ counts the number of ways to add the boxes that make up the Young diagram of $\mu$ to the diagram of $\nu$ in order to obtain that of
$\lambda$ subject to certain constraints (see [F-H, p. $455-56]$ ), and so each of these sums is actually finite.

The formula may be stated as

$$
\begin{equation*}
\chi_{S p_{6}}(\lambda) \chi_{S p_{6}}(\mu)=\pi_{S p_{6}}\left(\sum_{\zeta} \chi_{S p}((\lambda / \zeta) \cdot(\mu / \zeta))\right), \tag{7.4.2}
\end{equation*}
$$

where $\chi_{S p}(\nu)$ is a certain "universal character," which is defined for all partitions $\nu$ (while $\chi_{S p_{6}}$ is defined only for those with $\leq 3$ parts) and extended to finite formal sums by $\mathbb{Z}$ linearity, and $\pi_{S p_{6}}$ is a projection from the ring of universal characters to that of characters of $S p_{6}$. In our case, by the interpretation of the Littlewood-Richardson coefficients in terms of adding boxes to the Young diagram, we get

$$
\chi_{S p_{6}}(k) \chi_{S p_{6}}\left(\mu_{1} \mu_{2}^{2}\right)=\sum_{\nu, \kappa} \pi_{S p_{6}}\left(\chi_{S p}(\nu)\right),
$$

where the sum is over partitions $\nu$ and $\kappa$ satisfying

$$
0=\kappa_{4} \leq \nu_{4} \leq \kappa_{3} \leq \nu_{3} \leq \mu_{2}=\kappa_{2} \leq \nu_{2} \leq \kappa_{1} \leq \nu_{1}, \mu_{1}
$$

and

$$
\mu_{1}+\mu_{2}+\nu_{1}+\nu_{2}+\nu_{3}+\nu_{4}+\nu_{5}-2 \kappa_{1}-\kappa_{3}=k
$$

For $\nu$ having 3 or fewer parts, we have

$$
\pi_{S p_{6}}\left(\chi_{S p}(\nu)\right)=\chi_{S p_{6}}(\nu)
$$

For $\nu$ having exactly 4 parts, it is zero. So we find that (7.4.1) is equal to

$$
\begin{equation*}
(1-x y) \sum_{0 \leq \nu_{3} \leq \nu_{2} \leq \nu_{1}} \chi_{S p_{6}}(\nu) \sum_{\kappa_{3}=0}^{\nu_{3}} \sum_{\mu_{2}=\nu_{3}}^{\nu_{2}} \sum_{\kappa_{1}=\nu_{2}}^{\nu_{1}} \sum_{\mu_{1}=\kappa_{1}}^{\infty} x^{\mu_{1}+\nu_{1}+\nu_{2}+\nu_{3}-2 \kappa_{1}-2 \kappa_{3}}(-1)^{\mu_{2}+\nu_{1}+\nu_{2}+\nu_{3}} y^{\mu_{1}} . \tag{7.4.3}
\end{equation*}
$$

Now,

$$
\sum_{\mu_{1}=\kappa_{1}}^{\infty}(x y)^{\mu_{1}}=(x y)^{\kappa_{1}}(1-x y)^{-1}
$$

and

$$
\sum_{\mu_{2}=\nu_{3}}^{\nu_{2}}(-1)^{\mu_{2}}=(-1)^{\nu_{3}} \cdot \epsilon_{2}\left(\nu_{2}-\nu_{3}+1\right)
$$

Plugging these in, we obtain

$$
\begin{equation*}
\sum_{0 \leq \nu_{3} \leq \nu_{2} \leq \nu_{1}} \chi_{S p_{6}}(\nu) \epsilon_{2}\left(\nu_{2}-\nu_{3}+1\right) \sum_{\kappa_{3}=0}^{\nu_{3}} \sum_{\kappa_{1}=\nu_{2}}^{\nu_{1}} x^{\nu_{1}+\nu_{2}+\nu_{3}-\kappa_{1}-2 \kappa_{3}} y^{\kappa_{1}}(-1)^{\nu_{1}+\nu_{2}} . \tag{7.4.4}
\end{equation*}
$$

Summing over $\kappa_{1}$ yields

$$
\left(x^{-1} y\right)^{\nu_{2}} \frac{1-\left(x^{-1} y\right)^{\nu_{1}-\nu_{2}+1}}{1-x^{-1} y}
$$

We then write $\nu_{3}-2 \kappa_{3}=-\nu_{3}+2 r_{3}$, where $r_{3}=\nu_{3}-\kappa_{3}$, and sum $r_{3}$ from 0 to $\nu_{3}$ obtaining

$$
\frac{1-x^{2\left(\nu_{3}+1\right)}}{1-x^{2}}
$$

We now have

$$
\sum_{0 \leq \nu_{3} \leq \nu_{2} \leq \nu_{1}} \chi_{S p_{6}}(\nu) \epsilon_{2}\left(\nu_{2}-\nu_{3}+1\right) x^{\nu_{1}-\nu_{3}} y^{\nu_{2}}(-1)^{\nu_{1}-\nu_{2}} \frac{1-x^{2\left(\nu_{3}+1\right)}}{1-x^{2}} \frac{1-\left(x^{-1} y\right)^{\nu_{1}-\nu_{2}+1}}{1-x^{-1} y} .
$$

Translating this back to the other notation via the correspondence

$$
\chi_{S p_{6}}\left(\nu_{1}, \nu_{2}, \nu_{3}\right)=\operatorname{Tr}\left(\cdot \mid V_{\nu_{1}-\nu_{2}, \nu_{2}-\nu_{3}, \nu_{3}}^{S p_{6}}\right)
$$

and plugging in $\overline{s_{\pi}}$, we obtain the desired result. Lemma 7.1.1 is proved.

## 8. Unramified Computation: $E=F^{3}$

Finally, we come to the split case. This case is the most complicated of the three, because it involves decomposing the tensor product of 3 representations of $\operatorname{Spin}_{8}(\mathbb{C})$.
8.1. The identity. Let us set

$$
x_{i}=q^{-s_{i}} \quad \text { for } i=1,2 \text { or } 3 .
$$

Then the LHS of the identity in Theorem 5.5.3 is:

$$
\begin{aligned}
& \left(1-x_{1} x_{2} x_{3}\right)^{-1} \cdot \sum_{k_{i}=0}^{\infty}\left(k_{0}+1\right) \cdot V_{k_{0}, k_{1}, k_{2}, k_{3}} x_{1}^{k_{0}+k_{1}} x_{2}^{k_{0}+k_{2}} x_{3}^{k_{0}+k_{3}} \times \\
& \frac{1-\left(x_{1}^{-1} x_{2} x_{3}\right)^{k_{1}+1}}{1-\left(x_{1}^{-1} x_{2} x_{3}\right)} \frac{1-\left(x_{2}^{-1} x_{1} x_{3}\right)^{k_{2}+1}}{1-\left(x_{2}^{-1} x_{1} x_{3}\right)} \frac{1-\left(x_{3}^{-1} x_{1} x_{2}\right)^{k_{3}+1}}{1-\left(x_{3}^{-1} x_{1} x_{2}\right)}
\end{aligned}
$$

On the other hand, the RHS of the identity is:

$$
\prod_{i=1}^{3} \frac{1-x_{i}^{2}}{\operatorname{det}\left(1-x_{i} s_{\pi} \mid V_{\mu_{i}}\right)}
$$

Using the fact that (cf. [Br, Pg. 13])

$$
\operatorname{Sym}^{n}\left(V_{\mu_{i}}\right)=\operatorname{Sym}^{n-2}\left(V_{\mu_{i}}\right) \oplus V_{n \mu_{i}},
$$

this is equal to

$$
\sum_{\ell_{i}=0}^{\infty} V_{0, \ell_{1}, 0,0} \otimes V_{0,0, \ell_{2}, 0} \otimes V_{0,0,0, \ell_{3}} x_{1}^{\ell_{1}} x_{2}^{\ell_{2}} x_{3}^{\ell_{3}}
$$

Thus the identity we wish to prove is:

$$
\begin{gathered}
\sum_{\ell_{i}=0}^{\infty} V_{0, \ell_{1}, 0,0} \otimes V_{0,0, \ell_{2}, 0} \otimes V_{0,0,0, \ell_{3}} x_{1}^{\ell_{1}} x_{2}^{\ell_{2}} x_{3}^{\ell_{3}} \\
=\frac{1}{1-x_{1} x_{2} x_{3}} \sum_{k_{i}=0}^{\infty}\left(k_{0}+1\right) V_{k_{0}, k_{1}, k_{2}, k_{3}} x_{1}^{k_{0}+k_{1}} x_{2}^{k_{0}+k_{2}} x_{3}^{k_{0}+k_{3}} \\
\frac{1-\left(x_{1}^{-1} x_{2} x_{3}\right)^{k_{1}+1}}{1-\left(x_{1}^{-1} x_{2} x_{3}\right)} \frac{1-\left(x_{2}^{-1} x_{1} x_{3}\right)^{k_{2}+1}}{1-\left(x_{2}^{-1} x_{1} x_{3}\right)} \frac{1-\left(x_{3}^{-1} x_{1} x_{2}\right)^{k_{3}+1}}{1-\left(x_{3}^{-1} x_{1} x_{2}\right)} .
\end{gathered}
$$

8.2. Decomposition of tensor products. We make use of the fact (cf [GH, Lemma 3.3]) that

$$
V_{0,0, \ell_{2}, 0} \otimes V_{0,0,0, \ell_{3}}=\bigoplus_{k=0}^{\min \left(\ell_{2}, \ell_{3}\right)} V_{0, k, \ell_{2}-k, \ell_{3}-k}
$$

to rewrite the identity as

$$
\begin{gather*}
\sum_{\ell, m_{i}=0}^{\infty} V_{0, \ell, 0,0} \otimes V_{0, m_{1}, m_{2}, m_{3}} x_{1}^{\ell} x_{2}^{m_{1}+m_{2}} x_{3}^{m_{1}+m_{3}} \\
=\frac{1}{1-x_{1} x_{2} x_{3}} \sum_{k_{i}=0}^{\infty}\left(k_{0}+1\right) V_{k_{0}, k_{1}, k_{2}, k_{3}} x_{1}^{k_{0}+k_{1}} x_{2}^{k_{0}+k_{2}} x_{3}^{k_{0}+k_{3}} \\
\frac{1-\left(x_{1}^{-1} x_{2} x_{3}\right)^{k_{1}+1}}{1-\left(x_{1}^{-1} x_{2} x_{3}\right)} \frac{1-\left(x_{2}^{-1} x_{1} x_{3}\right)^{k_{2}+1}}{1-\left(x_{2}^{-1} x_{1} x_{3}\right)} \frac{1-\left(x_{3}^{-1} x_{1} x_{2}\right)^{k_{3}+1}}{1-\left(x_{3}^{-1} x_{1} x_{2}\right)} . \tag{8.2.1}
\end{gather*}
$$

We will be using a method of evaluating tensor products which is due to Black, King, and Wybourne [B-K-W]. They make use of an identification of certain weights of Spin ${ }_{8}$ with partitions. Specifically, we identify the highest weight of $V_{k_{0}, k_{1}, k_{2}, k_{3}}$ with the quadruple $\left(k_{0}+k_{1}+\frac{k_{2}+k_{3}}{2}, k_{0}+\frac{k_{2}+k_{3}}{2}, \frac{k_{2}+k_{3}}{2}, \frac{k_{3}-k_{2}}{2}\right)$ which, if $k_{3}-k_{2}$ is even and non-negative, is then identified with a partition by dropping any terminal zeros. We could, of course, have made such an identification in five other ways, and picking this one in particular privileges the eight dimensional representation $V_{0,1,0,0}$ over the other two, and fixes one particular identification of the representation ring of $\mathrm{SO}_{8}$ with a subring of that of $\mathrm{Spin}_{8}$. For $\mu$ a dominant weight, let $\chi_{\text {Spin }_{8}}(\mu)$ denote the character of the representation with highest weight $\mu$. If $\mu$ consists of integers, then this representation factors through $V_{0,1,0,0}$ and we may also write $\chi_{\mathrm{SO}_{8}}$.

Let $\iota$ denote the involution of Spin $_{8}$ which reverses the last two fundamental weights. Then the set of dominant weights is the union of the sets $\left\{\mu, \mu+\left(\frac{1}{2}\right)^{4}, \mu^{\iota},\left(\mu+\left(\frac{1}{2}\right)^{4}\right)^{\iota}\right\}$, where $\mu$ ranges over partitions. The union is disjoint, but if $\mu$ has fewer than 4 parts, then $\mu=\mu^{l}$. Hence the character of the LHS of (8.2.1) is given by

$$
\begin{align*}
& \sum_{\ell, \mu} \chi_{S O_{8}}(\ell) x_{1}^{\ell}\left(\chi_{S O_{8}}(\mu) x_{2}^{\mu_{1}-\mu_{4}} x_{3}^{\mu_{1}+\mu_{4}}+\chi_{S O_{8}}\left(\mu^{\iota}\right) x_{2}^{\mu_{1}+\mu_{4}} x_{3}^{\mu_{1}-\mu_{4}}+\right. \\
& \left.\chi_{S p i n_{8}}\left(\mu+\left(\frac{1}{2}\right)^{4}\right) x_{2}^{\mu_{1}-\mu_{4}} x_{3}^{\mu_{1}+\mu_{4}+1}+\chi_{S p i n_{8}}\left(\left(\mu+\left(\frac{1}{2}\right)^{4}\right)^{\iota}\right) x_{2}^{\mu_{1}+\mu_{4}+1} x_{3}^{\mu_{1}-\mu_{4}}\right)  \tag{8.2.2}\\
& \quad-\sum_{\ell, \mu: \mu_{4}=0} \chi_{S O_{8}}(\ell) \chi_{S O_{8}}(\mu) x_{1}^{\ell}\left(x_{2} x_{3}\right)^{\mu_{1}},
\end{align*}
$$

where $\ell$ is summed over $\mathbb{Z}_{\geq 0}$ and $\mu$ is summed over partitions satisfying $\mu_{2}=\mu_{3}$.
We describe the method of computing the necessary products of characters in brief here, and more completely in the appendix. First, the computation of the product of characters of $S_{\text {pin }}^{8}$ is reduced to one of characters of $G L_{4}$, which are indexed by pairs of partitions having at most four total parts. Such a pair of partitions also gives a character of $G L_{n}$ for each $n>4$. Each character is expressed as a determinant involving elementary symmetric polynomials and an identity in these determinantal expressions gives the product for all $n$ sufficiently large or, as we prefer, in a suitably defined projective limit. We thus obtain a finite sum of "universal characters" $\chi_{G L}(\bar{\tau} ; \nu)$. (We follow [B-K-W] in writing $(\bar{\tau} ; \nu)$, rather
than $(\tau, \nu))$. We then must analyze the contribution of each of these universal characters to the original product under projection.

We now write down a formula for the sum of universal characters that must be considered (cf. formula (11.4.4) in the appendix). Since one of our highest weights is the one-part partition $\ell$, the formula is a bit simpler than in the general case. It is:

$$
\begin{equation*}
\left.\sum_{\eta, \zeta} \chi_{G L}(\bar{\eta} ;(\ell /(\eta \cdot \zeta)) \cdot \mu)\right) . \tag{8.2.3}
\end{equation*}
$$

Here $\zeta$ and $\eta$ are summed over all partitions, but we only get contributions when each is a single integer and their sum does not exceed $\ell$. Making further use of the interpretation of the Littlewood-Richardson coefficients in terms of adding boxes to the Young diagram, we obtain

$$
\sum_{\nu, \kappa} \chi_{G L}\left(\overline{\ell^{\prime}(\mu, \nu, \kappa, \ell)} ; \nu\right),
$$

where the sum is over $\nu$ and $\kappa$ partitions satisfying

$$
\nu_{5} \leq \kappa_{4} \leq \mu_{4}, \nu_{4} \leq \kappa_{3} \leq \nu_{3} \leq \mu_{3}=\kappa_{2}=\mu_{2} \leq \nu_{2} \leq \kappa_{1} \leq \mu_{1}, \nu_{1},
$$

and subject to the condition that

$$
\ell^{\prime}(\mu, \nu, \kappa, \ell):=\ell+2 \kappa_{1}+2 \kappa_{3}+2 \kappa_{4}-\mu_{1}-\mu_{4}-\nu_{1}-\nu_{2}-\nu_{3}-\nu_{4}-\nu_{5}
$$

is non-negative. The precise manner in which a term $\chi_{G L}\left(\ell^{\prime} ; \nu\right)$ contributes to the product depends on the precise relationship between the highest weight of our second representation and the partition $\mu$. This is the topic of the next subsection. For now, we prove a formula for

$$
\begin{equation*}
\sum_{\mu, \ell} x^{\ell} y^{\mu_{1}-\mu_{4}} z^{\mu_{1}+\mu_{4}} \sum_{\nu, \kappa} \chi_{G L}\left(\overline{\left(\overline{\ell^{\prime}}(\mu, \nu, \kappa, \ell)\right.} ; \nu\right), \tag{8.2.4}
\end{equation*}
$$

where the sum in $\nu$ and $\kappa$ is as above. Let

$$
D\left(n_{3}, n_{4} ; x, z\right)=\left(1-x^{2 n_{3}}\right)\left(1-\left(x z^{-1}\right)^{n_{4}}\right)-(x z)^{n_{3}}\left(1-x^{2 n_{4}}\right)\left(1-\left(x z^{-1}\right)^{n_{3}}\right) .
$$

Then we prove

Proposition 8.2.5. The sum (8.2.4) is equal to

$$
F(x, y, z)^{-1} \sum_{\ell^{\prime}, \nu} \chi_{G L}\left(\bar{\ell}^{\prime} ; \nu\right) x^{x^{\prime}} c(\nu ; x, y, z)
$$

with

$$
F(x, y, z)=(1-x y)\left(1-x^{-1} y\right)(1-x z)\left(1-x z^{-1}\right)\left(1-x^{2}\right)
$$

and
$c(\nu ; x, y, z)=x^{\nu_{1}-\nu_{3}+\nu_{5}} y^{\nu_{2}} z^{\nu_{4}}\left(\nu_{2}-\nu_{3}+1\right)\left(1-\left(x^{-1} y\right)^{\nu_{1}-\nu_{2}+1}\right) D\left(\nu_{3}-\nu_{4}+1, \nu_{4}-\nu_{5}+1 ; x, z\right)$.
Proof. We first bring the sum over $\ell$ inside, and get

$$
\sum_{\nu, \ell^{\prime}} \chi_{G L}\left(\overline{\ell^{\prime}} ; \nu\right) x^{\ell^{\prime}} \sum_{\mu} y^{\mu_{1}} z^{\mu_{4}} P\left(\mu, \nu ; x_{1}\right),
$$

where

$$
P(\mu, \nu ; x)=\sum_{\kappa} x^{\mu_{1}+\mu_{4}+\nu_{1}+\nu_{2}+\nu_{3}+\nu_{4}+\nu_{5}-2 \kappa_{1}-2 \kappa_{3}-2 \kappa_{4}} .
$$

Observe that

$$
P(\mu, \nu ; x)=P_{1}\left(\mu_{1}, \nu_{1}, \nu_{2} ; x\right) P_{2}\left(\mu_{4}, \nu_{3}, \nu_{4}, \nu_{5} ; x\right)
$$

where

$$
P_{1}\left(\mu_{1}, \nu_{1}, \nu_{2} ; x\right)=\sum_{\kappa_{1}=\nu_{2}}^{\min \left(\mu_{1}, \nu_{1}\right)} x^{\mu_{1}+\nu_{1}+\nu_{2}-2 \kappa_{1}}
$$

and

$$
P_{2}\left(\mu_{4}, \nu_{3}, \nu_{4}, \nu_{5} ; x\right)=\sum_{\kappa_{4}=\nu_{5}}^{\min \left(\mu_{4}, \nu_{4}\right)} \sum_{\kappa_{3}=\max \left(\mu_{4}, \nu_{4}\right)}^{\nu_{3}} x^{\mu_{4}+\nu_{3}+\nu_{4}+\nu_{5}-2 \kappa_{4}-2 \kappa_{3}}
$$

Now, let

$$
Q_{1}\left(\nu_{1}, \nu_{2} ; x, y\right)=\sum_{\mu_{1}=\nu_{2}}^{\infty} P_{1}\left(\mu_{1}, \nu_{1}, \nu_{2} ; x\right) y^{\mu_{1}}=\sum_{\kappa_{1}=\nu_{2}}^{\nu_{1}} \sum_{\mu_{1}=\kappa_{1}}^{\infty} x^{\nu_{1}+\nu_{2}+\mu_{1}-2 \kappa_{1}} y^{\mu_{1}}
$$

and let

$$
Q_{2}\left(\nu_{3}, \nu_{4}, \nu_{5} ; x, z\right)=\sum_{\mu_{4}=\nu_{5}}^{\nu_{3}} P_{2}\left(\mu_{4}, \nu_{3}, \nu_{4}, \nu_{5} ; x\right) z^{\mu_{4}}=\sum_{\kappa_{4}=\nu_{5}}^{\nu_{4}} \sum_{\kappa_{3}=\nu_{4}}^{\nu_{3}} \sum_{\mu_{4}=\kappa_{4}}^{\kappa_{3}} x^{\mu_{4}+\nu_{3}+\nu_{4}+\nu_{5}-2 \kappa_{4}-2 \kappa_{3}} z^{\mu_{4}}
$$

Then we have

$$
c(\nu ; x, y, z)=\left(\nu_{2}-\nu_{3}+1\right) Q_{1}\left(\nu_{1}, \nu_{2} ; x, y\right) Q_{2}\left(\nu_{3}, \nu_{4}, \nu_{5} ; x, z\right)
$$

The $\left(\nu_{2}-\nu_{3}+1\right)$ comes because in original sum we had a free variable (say, $\mu_{3}$ ) ranging from $\nu_{3}$ to $\nu_{2}$. The $Q_{i}$ can be evaluated explicitly:

$$
Q_{1}\left(\nu_{1}, \nu_{2} ; x, y\right)=(1-x y)^{-1}\left(1-x^{-1} y\right)^{-1} x^{\nu_{1}} y^{\nu_{2}}\left(1-\left(x^{-1} y\right)^{\nu_{1}-\nu_{2}+1}\right)
$$

and

$$
\begin{gathered}
Q_{2}\left(\nu_{3}, \nu_{4}, \nu_{5} ; x, z\right)=(1-x z)^{-1}\left(1-x z^{-1}\right)^{-1}\left(1-x^{2}\right)^{-1} z^{\nu_{4}} x^{-\nu_{3}+\nu_{5}} \\
\left(\left(1-x^{2\left(\nu_{3}-\nu_{4}+1\right)}\right)\left(1-\left(x z^{-1}\right)^{\nu_{4}-\nu_{5}+1}\right)-(x z)^{\nu_{3}-\nu_{4}+1}\left(1-x^{2\left(\nu_{4}-\nu_{5}+1\right)}\right)\left(1-\left(x z^{-1}\right)^{\nu_{3}-\nu_{4}+1}\right)\right)
\end{gathered}
$$

The result follows.

In the next subsection, we describe how each $\chi_{G L}\left(\overline{\ell^{\prime}} ; \nu\right)$ gives a contribution to the product

$$
\chi_{S O_{8}}(\ell) \chi_{\operatorname{Spin}_{8}}(\xi), \quad \text { with } \xi=\left(\mu+\varepsilon^{4}\right) \text { or }\left(\mu+\varepsilon^{4}\right)^{\iota} \text { and } \varepsilon=0 \text { or } \frac{1}{2}
$$

Just as the precise contribution depends on $\varepsilon$ and the presence or absence of $\iota$, so too does the polynomial which will accompany it in (8.2.2). As may be seen from (8.2.2), and the fact that $F(x, y, z)=F\left(x, y, z^{-1}\right)$, this will be

$$
F\left(x_{1}, x_{2} x_{3}, x_{2} x_{3}^{-1}\right)^{-1} x_{1}^{\ell^{\prime}} \times\left\{\begin{array}{l}
c\left(\nu ; x_{1}, x_{2} x_{3}, x_{2}^{-1} x_{3}\right) \text { for } \mu \\
c\left(\nu ; x_{1}, x_{2} x_{3}, x_{2} x_{3}^{-1}\right) \text { for } \mu^{\iota} \\
x_{3} c\left(\nu ; x_{1}, x_{2} x_{3}, x_{2}^{-1} x_{3}\right) \text { for } \mu+\left(\frac{1}{2}\right) \\
x_{2} c\left(\nu ; x_{1}, x_{2} x_{3}, x_{2} x_{3}^{-1}\right) \text { for }\left(\mu+\left(\frac{1}{2}\right)\right)^{\iota}
\end{array}\right.
$$

We also record for later use the analogous statement for the smaller sum being subtracted at the end of (8.2.2). The proof is easily obtained by adapting that of the previous proposition:

Lemma 8.2.6. We have

$$
\sum_{\ell, \mu: \mu_{4}=0} x_{1}^{\ell}\left(x_{2} x_{3}\right)^{\mu_{1}} \chi_{G L}\left(\overline{\ell^{\prime}(\mu, \nu, \kappa, \ell)} ; \nu\right)=\sum_{\ell^{\prime}, \nu} \chi_{G L}\left(\bar{\ell}^{\prime} ; \nu\right) x_{1}^{\ell^{\prime}} c^{\prime}\left(\nu ; x_{1}, x_{2} x_{3}\right)
$$

where $\nu$ can have at most four parts, and

$$
c^{\prime}\left(\nu ; x_{1}, x_{2} x_{3}\right)=\frac{\left(\nu_{2}-\nu_{3}+1\right) x_{1}^{\nu_{1}-\nu_{3}+\nu_{4}}\left(x_{2} x_{3}\right)^{\nu_{2}}\left(1-\left(x_{1}^{-1} x_{2} x_{3}\right)^{\nu_{1}-\nu_{2}+1}\right)\left(1-x_{1}^{2\left(\nu_{3}-\nu_{4}+1\right)}\right)}{\left(1-x_{1} x_{2} x_{3}\right)\left(1-x_{1}^{-1} x_{2} x_{3}\right)\left(1-x_{1}^{2}\right)} .
$$

8.3. From $\chi_{G L}$ 's to $\chi_{S p i n_{8}}$ 's. Now, we briefly describe the contribution of $\chi_{G L}(\bar{\ell} ' ; \nu)$ to $\chi_{S O_{8}}(\ell) \chi_{\text {Spin }_{8}}\left(\mu+\varepsilon^{4}\right)$. The answer is always zero or $\pm_{\text {Spin }_{8}}(\xi)$ for some weight $\xi$. The corresponding contribution to $\chi_{S O_{8}}(\ell) \chi_{\text {Spin }}^{8}\left(\left(\mu+\varepsilon^{4}\right)^{\iota}\right)$ would then be zero or $\pm_{\text {Spin }_{8}}\left(\xi^{\iota}\right)$. First, there is a projection from universal characters to characters of $G L_{4}$ which sends $\chi_{G L}(\bar{\tau} ; \nu)$ either to 0 or to $\pm \chi_{G L_{4}}\left(\bar{\tau}^{\circ} ; \nu^{\circ}\right)$ for a pair of partitions $\left(\bar{\tau}^{\circ} ; \nu^{\circ}\right)$ with at most four total parts. In our case, $\tau=\ell^{\prime}$, and $\nu$ can have at most five parts. Among such pairs, the only ones which have more than four total parts, (recall that $\ell^{\prime}=0$ gives the empty partition having no parts) and are not killed are those of the form ( $\overline{1} ; \sigma 1$ ), where $\sigma 1$ is the partition obtained by appending a 1 to the partition $\sigma$ having exactly four parts. The projection of $\chi_{G L}(\overline{1} ; \sigma 1)$ is $-\chi_{G L_{4}}(0 ; \sigma)$. If the total number of parts is already less than four, then $\chi_{G L}\left(\bar{\ell}^{\prime} ; \nu\right)$ projects to $\chi_{G L_{4}}\left(\bar{\ell}^{\prime} ; \nu\right)$.

The parametrization of representations by pairs of partitions is such that $\chi_{G L_{4}}\left(\bar{\ell}^{\prime} ; \nu\right)$ is the character of the representation whose highest weight $\omega$ is $\nu$ if $\ell^{\prime}=0$ and $\left(\nu_{1}, \nu_{2}, \nu_{3},-\ell^{\prime}\right)$ otherwise. (Here we append terminal zeros to $\nu$ if it has fewer than 3 parts.) Now, we must add $\varepsilon^{4}$ to $\omega$. The contribution to $\chi_{S O_{8}}(\ell) \chi_{S p i n}\left(\mu+\varepsilon^{4}\right)$ is as follows:

$$
\left\{\begin{array}{l}
\chi_{\text {Spin }_{8}\left(\nu+\varepsilon^{4}\right), \text { if } \ell^{\prime}=0}^{0, \text { if } \omega+\varepsilon^{4}+\rho \text { has a nontrivial stabilizer in the Weyl group of } \text { Spin }_{8},}  \tag{8.3.1}\\
\operatorname{sgn}(w) \chi_{\text {Spin }}(\xi), \text { if } \omega+\varepsilon^{4}+\rho=w(\xi+\rho) \text { for } \xi \text { dominant. }
\end{array}\right.
$$

Here $\rho$ is half the sum of the positive roots of $\operatorname{Spin}_{8}$, which corresponds to the quadruple $(3,2,1,0)$. The last case will produce a weight of the form $\left(\sigma+\varepsilon^{4}\right)^{\iota}$, with $\sigma$ a partition. Indeed, $\left(\sigma+\varepsilon^{4}\right)^{\iota}$ arises directly as the weight corresponding to $\left(\overline{\sigma_{4}+2 \varepsilon} ; \sigma_{1}, \sigma_{2}, \sigma_{3}\right)$ and the weights corresponding to the three other pairs:
$\left(\overline{\sigma_{3}+1+2 \varepsilon} ; \sigma_{1}, \sigma_{2}, \sigma_{4}-1\right),\left(\overline{\sigma_{2}+2+2 \varepsilon} ; \sigma_{1}, \sigma_{3}-1, \sigma_{4}-1\right),\left(\overline{\sigma_{1}+3+2 \varepsilon} ; \sigma_{2}-1, \sigma_{3}-1, \sigma_{4}-1\right)$
are related to it by Weyl elements $w$ with signs,-+ and - respectively, provided $\sigma$ has four parts. Hence these terms in the sum for $\left(\mu+\varepsilon^{4}\right)^{\iota}$ contribute to $\left(\sigma+\varepsilon^{4}\right)$.
8.4. Completion of proof. We now fix a partition $\sigma$ and compute the coefficient of $\chi_{\text {Spin }}^{8}(\sigma+$ $\left.\varepsilon^{4}\right)$ in (8.2.2). Since $\chi_{\text {Spin }}^{8}\left(\sigma+\varepsilon^{4}\right)$ is the character of $V_{\sigma_{2}-\sigma_{3}, \sigma_{1}-\sigma_{2}, \sigma_{3}-\sigma_{4}, \sigma_{3}+\sigma_{4}+2 \varepsilon}$, we should get

$$
\begin{gathered}
F\left(x_{1}, x_{2} x_{3}, x_{2} x_{3}^{-1}\right)^{-1}\left(1-x_{1}^{2}\right)\left(\sigma_{2}-\sigma_{3}+1\right) x_{1}^{\sigma_{1}-\sigma_{3}} \sigma_{2}^{\sigma_{2}-\sigma_{4}} x_{3}^{\sigma_{2}+\sigma_{4}+2 \varepsilon} \\
\left(1-\left(x_{1}^{-1} x_{2} x_{3}\right)^{\sigma_{1}-\sigma_{2}+1}\right)\left(1-\left(x_{1} x_{2}^{-1} x_{3}\right)^{\sigma_{3}-\sigma_{4}+1}\right)\left(1-\left(x_{1} x_{2} x_{3}^{-1}\right)^{\sigma_{3}+\sigma_{4}+2 \varepsilon+1}\right) .
\end{gathered}
$$

Furthermore, once we check this, we are done, since both sides of (8.2.1) are symmetric if we replace the subscript 2 everywhere by 3 and vice versa.

Clearing denominators, we just need to check some identities of polynomials. As will be evident from the discussion above, there are fewer terms if $\sigma$ has fewer than four parts. However, if we append zeros to $\sigma$, the case $\sigma_{4}=0$ is recovered from the general case. For example, the term corresponding to ( $\overline{\sigma_{3}+1+2 \varepsilon} ; \sigma_{1}, \sigma_{2}, \sigma_{4}-1$ ) should not be there, but there
is no harm in including it, because $D\left(n_{3}, 0 ; x, z\right)$ is equal to zero anyway. Thus, there are only two identities to check.

First, we perform a simplification that is relevant to both; namely one shows that

$$
\begin{gathered}
x_{1}^{\sigma_{1}-\sigma_{3}+\sigma_{4}} x_{2}^{\sigma_{2}} x_{3}^{\sigma_{2}}\left(\sigma_{2}-\sigma_{3}+1\right)\left(1-\left(x_{1}^{-1} x_{2} x_{3}\right)^{\sigma_{1}-\sigma_{2}+1}\right) D\left(\sigma_{3}+1,1 ; x_{1}, x_{2} x_{3}^{-1}\right) \\
\quad-x_{1}^{\sigma_{1}+\sigma_{3}-\sigma_{4}+2} x_{2}^{\sigma_{2}} x_{3}^{\sigma_{2}}\left(\sigma_{2}-\sigma_{4}+2\right)\left(1-\left(x_{1}^{-1} x_{2} x_{3}\right)^{\sigma_{1}-\sigma_{2}+1}\right) D\left(\sigma_{4}, 1 ; x_{1}, x_{2} x_{3}^{-1}\right) \\
+ \\
x_{1}^{\sigma_{1}+\sigma_{2}-\sigma_{4}+3} x_{2}^{\sigma_{3}-1} x_{3}^{\sigma_{3}-1}\left(\sigma_{3}-\sigma_{4}+1\right)\left(1-\left(x_{1}^{-1} x_{2} x_{3}\right)^{\sigma_{1}-\sigma_{3}+2}\right) D\left(\sigma_{4}, 1 ; x_{1}, x_{2} x_{3}^{-1}\right) \\
- \\
x_{1}^{\sigma_{1}+\sigma_{2}-\sigma_{4}+3} x_{2}^{\sigma_{3}-1} x_{3}^{\sigma_{3}-1}\left(\sigma_{3}-\sigma_{4}+1\right)\left(1-\left(x_{1}^{-1} x_{2} x_{3}\right)^{\sigma_{2}-\sigma_{3}+1}\right) D\left(\sigma_{4}, 1 ; x_{1}, x_{2} x_{3}^{-1}\right)
\end{gathered}
$$

is equal to

$$
\begin{gathered}
\left(\sigma_{2}-\sigma_{3}+1\right) x_{1}^{\sigma_{1}-\sigma_{3}+\sigma_{4}}\left(x_{2} x_{3}\right)^{\sigma_{2}}\left(1-\left(x_{1}^{-1} x_{2} x_{3}\right)^{\sigma_{1}-\sigma_{2}+1}\right) \\
\left(D\left(\sigma_{3}+1,1 ; x_{1}, x_{2} x_{3}^{-1}\right)-x_{1}^{2\left(\sigma_{3}-\sigma_{4}+1\right)} D\left(\sigma_{4}, 1 ; x_{1}, x_{2} x_{3}^{-1}\right)\right)
\end{gathered}
$$

Canceling from both sides of (8.2.1) the expression

$$
\left(\sigma_{2}-\sigma_{3}+1\right) x_{1}^{\sigma_{1}-\sigma_{3}} x_{2}^{\sigma_{2}-\sigma_{4}} x_{3}^{\sigma_{2}+\sigma_{4}+2 \varepsilon}\left(1-\left(x_{1}^{-1} x_{2} x_{3}\right)^{\sigma_{1}-\sigma_{2}+1}\right),
$$

which appears in $c\left(\sigma ; x_{1}, x_{2} x_{3}, x_{2}^{-1} x_{3}\right), c\left(\sigma 1 ; x_{1}, x_{2} x_{3}, x_{2}^{-1} x_{3}\right)$ and $c^{\prime}\left(\sigma ; x_{1}, x_{2} x_{3}, x_{2} x_{3}^{-1}\right)$ (when $\varepsilon=0$ ), and making use of (8.3.1), the identities to be checked boil down to:
(i) corresponding to $\varepsilon=0$,

$$
\begin{gathered}
D\left(\sigma_{3}-\sigma_{4}+1, \sigma_{4}+1 ; x_{1}, x_{2}^{-1} x_{3}\right)-x_{1}^{2} D\left(\sigma_{3}-\sigma_{4}+1, \sigma_{4} ; x_{1}, x_{2}^{-1} x_{3}\right)+\left(x_{1} x_{2} x_{3}^{-1}\right)^{\sigma_{4}} \\
\left(D\left(\sigma_{3}+1,1 ; x_{1}, x_{2} x_{3}^{-1}\right)-x_{1}^{2\left(\sigma_{3}-\sigma_{4}+1\right)} D\left(\sigma_{4}, 1 ; x_{1}, x_{2} x_{3}^{-1}\right)-\left(1-x_{1} x_{2}^{-1} x_{3}\right)\left(1-x_{1} x_{2} x_{3}^{-1}\right)\left(1-x_{1}^{2\left(\sigma_{3}-\sigma_{4}+1\right)}\right)\right) \\
=\left(1-\left(x_{1} x_{2}^{-1} x_{3}\right)^{\sigma_{3}-\sigma_{4}+1}\right)\left(1-\left(x_{1} x_{2}^{-1} x_{3}^{-1}\right)^{\sigma_{3}+\sigma_{4}+1}\right)\left(1-x_{1}^{2}\right)
\end{gathered}
$$

(ii) corresponding to $\varepsilon=\frac{1}{2}$,

$$
\begin{gathered}
D\left(\sigma_{3}-\sigma_{4}+1, \sigma_{4}+1 ; x_{1}, x_{2}^{-1} x_{3}\right)-x_{1}^{2} D\left(\sigma_{3}-\sigma_{4}+1, \sigma_{4} ; x_{1}, x_{2}^{-1} x_{3}\right) \\
+\left(x_{1} x_{2} x_{3}^{-1}\right)^{\sigma_{4}+1}\left(D\left(\sigma_{3}+1,1 ; x_{1}, x_{2} x_{3}^{-1}\right)-x_{1}^{2\left(\sigma_{3}-\sigma_{4}+1\right)} D\left(\sigma_{4}, 1 ; x_{1}, x_{2} x_{3}^{-1}\right)\right) \\
=\left(1-\left(x_{1} x_{2}^{-1} x_{3}\right)^{\sigma_{3}-\sigma_{4}+1}\right)\left(1-\left(x_{1} x_{2}^{-1} x_{3}\right)^{\sigma_{3}+\sigma_{4}+2}\right)\left(1-x_{1}^{2}\right) .
\end{gathered}
$$

Each of these may be checked by hand.
This completes the proof of Theorem 5.5.3 in the split case.

## 9. Ramified Factors

Finally, we shall address the analytic properties of the local zeta factor at a ramified place $v \in S$. Thus, we continue to work locally and suppress $v$ from the notations.
9.1. The local zeta integral. Recall that:

$$
\begin{align*}
Z(\varphi, \Phi, f, \underline{s}) & =\int_{V_{E}^{\prime}(F) \backslash G_{E}(F)} W_{V_{E}^{\prime}, \psi}(g \cdot \varphi) \cdot \overline{L_{\chi_{0}}(g \cdot f)} \cdot \widehat{\Phi_{s}}(g) d g \\
& =\int_{V_{E}^{\prime}(F) \backslash G_{E}(F)} \int_{U_{-\alpha}(F)} W_{V_{E}^{\prime}, \psi}(g \cdot \varphi) \cdot \overline{L_{\chi_{0}}(g \cdot f)} \cdot \psi(u) \cdot \Phi_{s}(u g) d u d g  \tag{9.1.1}\\
& =\int_{S_{E}} \int_{K} \int_{U_{-\alpha}} W_{V_{E}^{\prime}, \psi}(t k \cdot \varphi) \cdot \overline{L_{\chi_{0}}(t k \cdot f)} \cdot \psi(u) \cdot \Phi_{s}(u t k) \cdot \delta_{B_{E}^{\prime}}(t)^{-1} d u d k d t . \tag{9.1.2}
\end{align*}
$$

In this section, by the local zeta integral $Z(\varphi, \Phi, f, \underline{s})$, we mean this last integral (9.1.2).
Assuming absolute convergence of (9.1.2), the manipulations above would be justified and we can collapse the two integrals in (9.1.1) to get:

$$
\begin{aligned}
& Z(\varphi, \Phi, f, \underline{s}) \\
= & \int_{N_{E}(F) \backslash G_{E}(F)} W_{V_{E}^{\prime}, \psi}(g \cdot \varphi) \cdot \overline{L_{\chi_{0}}(g \cdot f)} \cdot \Phi_{s}(g) d g \\
= & \int_{V_{E}(F) \backslash G_{E}(F)} \Phi_{\underline{s}}(g) \cdot\left(\int_{U_{\alpha}} W_{V_{E}^{\prime}, \psi}(u g \cdot \varphi) \cdot \overline{L_{\chi_{0}}(u g \cdot f)} d u\right) d g \\
= & \int_{K} \Phi(k)\left(\int_{S_{E}} \chi_{\underline{s}}(t) \cdot \delta_{B_{E}}(t)^{-1} \cdot \int_{U_{\alpha}} W_{V_{E}^{\prime}, \psi}(u t k \cdot \varphi) \cdot \overline{L_{\chi_{0}}(u t k \cdot f)} d u d t\right) d k .
\end{aligned}
$$

Setting $Y(\varphi, f, \underline{s}, k)$ to be the inner integral, we would have

$$
Z(\varphi, \Phi, f, \underline{s})=\int_{K} \Phi(k) \cdot Y(\varphi, f, \underline{s}, k) d k
$$

The above formal manipulation is justified by the following proposition, which is the main result of this section:

Proposition 9.1.3. (i) The local zeta integral $Z(\varphi, \Phi, f, \underline{s})$ (i.e. the integral (9.1.2)) converges absolutely for $\underline{s} \in \Omega_{R}$ (cf. (2.6)) when $R$ is sufficiently large. In particular, $Y(\varphi, f, \underline{s}, k)$ converges absolutely in $\Omega_{R}$ as well, and we have

$$
Z(\varphi, \Phi, f, \underline{s})=\int_{K} \Phi(k) \cdot Y(\varphi, f, \underline{s}, k) d k
$$

for $\underline{s} \in \Omega_{R}$.
(ii) Both $Y(\varphi, f, \underline{s},-)$ and $Z(\varphi, \Phi, f, \underline{s})$ admit meromorphic continuation to the whole of $\mathbb{C}^{r-1}$ where $r=\operatorname{rank}_{F}\left(G_{E}\right)$.
(iii) For a fixed $\underline{s_{0}} \in \mathbb{C}^{r-1}$, there is a choice of $\varphi, \Phi$ and $f$ such that $Z(\varphi, \Phi, f, \underline{s})$ is holomorphic at $\underline{s_{0}}$ and non-zero there.

As a consequence, we have:
Corollary 9.1.4. The partial Spin L-function $L^{S}\left(\underline{s}, \pi, V_{\text {Spin }, E}\right)$ admits meromorphic continuation to $\mathbb{C}^{r-1}$.

The rest of this section is devoted to the proof of the proposition.
9.2. Asymptotics of functions. With $r=\operatorname{rank}_{F}\left(G_{E}\right)$, we identify $S_{E}(F)$ with $\left(F^{\times}\right)^{r}$ using the fundamental coweights $\left\{\mu_{i}^{\prime}\right\}$ as described in (5.1). We first note the following lemmas:
Lemma 9.2.1. There are finitely many finite functions $\xi_{i}$ on $\left(F^{\times}\right)^{r}$ such that for any $\varphi \in \pi$, one can find Schwarz-Bruhat functions $\varphi_{i}$ (depending on $\varphi$ ) on $F^{r} \times K$ satisfying

$$
W_{V_{E}^{\prime}, \psi}(t k \cdot \varphi)=\sum_{i} \varphi_{i}(t, k) \cdot \xi_{i}(t)
$$

This lemma is well-known. See [JS, $\S 4$, Props. 1 and 3] and $[S, \S 4$, Thm. 1]. We have a similar lemma for the asymptotics of the linear functional $L_{\chi_{0}}$ on the minimal representation.

Lemma 9.2.2. There are finitely many finite functions $\eta_{i}$ on $\left(F^{\times}\right)^{r}$ such that given any $f \in \Pi_{E}$, one can find Schwarz-Bruhat functions $f_{i}$ on $F \times K$ satisfying

$$
L_{\chi_{0}}(t k \cdot f)=\sum_{i} f_{i}\left(t_{0}, k\right) \cdot \eta_{i}(t)
$$

Proof. This follows from the construction of the linear functional $L_{\chi_{0}}$ described in the proof of Prop. 3.4.1. Indeed, if $L_{E}^{\prime} U_{E}^{\prime}$ is the 3-step parabolic containing $B_{E}^{\prime}$, then $L_{\chi_{0}}$ was described as the composition of an $U_{E}^{\prime}$-invariant map to a representation $\sigma$ of $L_{E}^{\prime}$ followed by the Whittaker functional on $\sigma$. Thus the lemma follows by the analog of the previous lemma for Whittaker functions on $L_{E}$.
9.3. Proof of Prop. 9.1.3(i). We are now ready to prove Prop. 9.1.3(i), i.e. that the integral

$$
\int_{K} \int_{S_{E}} \int_{U_{-\alpha}}\left|W_{V_{E}^{\prime}, \psi}(t k \cdot \varphi)\right| \cdot\left|\overline{L_{\chi_{0}}(t k \cdot f)}\right| \cdot\left|\Phi_{s}(u t k)\right| \cdot \delta_{B_{E}^{\prime}}(t)^{-1} d u d t d k
$$

converges when $\underline{s} \in \Omega_{R}$ with $R \gg 0$.
If we commute $t$ across $u$ in the above integral and change variables in $u$, as we did in the calculation before Prop. 5.5.2, we see that we need to prove the convergence of

$$
\int_{K} \int_{S_{E}}\left|W_{V_{E}^{\prime}, \psi}(t k \cdot \varphi)\right| \cdot \mid \overline{L_{\chi_{0}}(t k \cdot f) \mid} \cdot \delta_{\underline{s}}(t)^{-1} \cdot\left(\int_{U_{-\alpha}}\left|\Phi_{\underline{s}}(u k)\right| d u\right) d t d k
$$

where $\delta_{\underline{s}}$ is an explicit character of $t$ which can be read off from the computation before Prop. 5.5.2.

Now the integral over $U_{-\alpha}$ is a standard intertwining operator which converges in $\Omega_{R}$ for $R \gg 0$ and defines a smooth function of $k$. On the other hand, the convergence of the integrals over $S_{E}$ and $K$ follows from Lemmas 9.2 .1 and 9.2.2. This proves Prop. 9.1.3(i).
9.4. Proof of Prop. 9.1.3(ii). To prove meromorphic continuation, we first note:

Lemma 9.4.1. (i) When $\underline{s} \in \Omega_{R}$ with $R \gg 0$, the integral defining $\widehat{\Phi_{\underline{s}}}(g)$ converges absolutely.
(ii) It admits a holomorphic continuation to $\mathbb{C}^{r-1}$.
(iii) Moreover, there are finitely many finite functions $\chi_{\underline{s}, i}$ on $\left(F^{\times}\right)^{r}$ depending holomorphically on $\underline{s}$ such that for any flat section $\Phi_{\underline{s}}$, one can fin $\bar{d}$ Schwartz-Bruhat functions $\phi_{, \underline{s}, i}$ on $\left(F^{\times}\right)^{r-1} \times K$ depending holomorphically on s. satisfying

$$
\widehat{\Phi}_{\underline{s}}(t)=\sum_{i} \phi_{\underline{s}, i}\left(t_{1}, \ldots, t_{r-1}, k\right) \cdot \chi_{\underline{s}, i}(t)
$$

The finite functions $\chi_{\underline{s}, i}$ are precisely those which appear in Proposition 5.5.2.

Proof. (i) The absolute convergence of the integral defining $\widehat{\Phi}_{\underline{s}}$ has been addressed in the proof of Prop. 9.1.3(i) above.
(ii) Observe that the value $\widehat{\Phi_{\underline{s}}}(g)$ is obtained by restricting $g \cdot \Phi_{\underline{s}}$ to $M_{E}$, thus obtaining a section in a family of principal series of $M_{E}$, followed by applying the Jacquet integral for this family of principal series (which is a Whittaker functional). The desired holomorphic continuation was shown in Jacquet's thesis [J1].
(iii) This follows by an analog of Lemma 9.2.1 with parameters. Such a result is proved in [S, §4, Thm. 4], as well as in the recent article [J2].

As a consequence of Lemmas 9.2.1, 9.2.2 and 9.4.1, the local zeta integral is equal (for $\underline{s} \in \Omega_{R}$ ) to a sum of various integrals over $S_{E} \times K$ of finite functions $\xi_{i}$ on $S_{E}$ against Schwarz-Bruhat functions $\phi_{i}$ on $F^{r} \times K$. If the Schwarz functions $\phi_{i}$ are independent of $\underline{s}$, the meromorphic continuation of this type of integrals is well-known (cf. [JS, §3]). Indeed, in the p-adic case, the integral is easily seen to be equal to a rational function in $q^{-s_{i}}$. For a general local field, it was shown in [JS, §3] that the resulting meromorphic function is the product of various abelian $L$-functions (which are independent of the $\phi_{i}$ 's) and an entire function depending on the $\phi_{i}$ 's.

Now in our case, the Schwarz functions $\phi_{i}$ do depend holomorphically on $\underline{s}$ (via Lemma 9.4.1). The meromorphic continuation of our integrals then follows from the case discussed in the previous paragraph and [J2, Lemma 1, Pg. 377]. This proves the meromorphic continuation of $Z(\varphi, \Phi, f, \underline{s})$. The same argument gives the meromorphic continuation of the integral $Y$; one simply omits the integration over $K$. This proves Prop. 9.1.3(ii).

### 9.5. Proof of Prop. 9.1.3(iii). Recall that

$$
Z(\varphi, \Phi, f, \underline{s})=\int_{K} \Phi(k) \cdot Y(\varphi, f, \underline{s}, k) d k
$$

The function $\Phi$ is an arbitrary smooth function on $K$ subject to the condition that

$$
\Phi(l k)=\Phi(k) \quad \text { for all } l \in K \cap Q_{E}(F) .
$$

The function $Y(\varphi, f, \underline{s}, k)$ on $K$ is easily seen to be left invariant under $K \cap U_{E}(F)$. Thus to show that data can be chosen to ensure the non-vanishing of $Z(\varphi, \Phi, f, \underline{s})$ at $\underline{s}=\underline{s_{0}}$, we need to show that the integral

$$
Y^{\prime}(\varphi, f, \underline{s})=\int_{K \cap L_{E}(F)} Y(\varphi, f, \underline{s}, l) d l
$$

is non-zero for some choices of $f$ and $\varphi$. Note that since $Y$ has meromorphic continuation to $\mathbb{C}^{r-1}$, so does the integral $Y^{\prime}$.

Now we are going to massage the expression for $Y^{\prime}$ as follows:

$$
\begin{aligned}
& Y^{\prime}(\varphi, f, \underline{s}) \\
= & \int_{K \cap L_{E}}\left(\int_{S_{E}} \int_{U_{\alpha}} \chi_{\underline{s}}(t) \cdot \delta_{B_{E}}(t)^{-1} \cdot W_{V_{E}^{\prime}, \psi}(u t k \cdot \varphi) \cdot \overline{L_{\chi_{0}}(u t k \cdot f)} d u d t\right) d k \\
= & \int_{U_{\beta} \backslash L_{E}} \int_{U_{\alpha}} \chi_{\underline{s}}(l) \cdot \delta_{Q_{E}}(l)^{-1} \cdot W_{V_{E}^{\prime}, \psi}(u l \cdot \varphi) \cdot \overline{L_{\chi_{0}}(u l \cdot f)} d u d l \\
= & \int_{T_{E}} \int_{U_{-\beta}} \int_{U_{\alpha}} \chi_{\underline{s}}(t) \cdot \delta_{Q_{E}}(t)^{-1} \cdot W_{V_{E}^{\prime}, \psi}\left(u_{\alpha} u_{-\beta} t \cdot \varphi\right) \cdot \overline{L_{\chi_{0}}\left(u_{\alpha} u_{-\beta} t \cdot f\right)} d u_{\alpha} d u_{-\beta} d t .
\end{aligned}
$$

In the above, to obtain the second equality, we have used the Iwasawa decomposition for $L_{E}$, and to obtain the last equality, we have used the Bruhat decomposition to replace the integral over $U_{\beta} \backslash L_{E}$ by the integral over the open dense subset $U_{-\beta} \cdot T_{E}$.

Now the above manipulations are initially valid for $\underline{s} \in \Omega_{R}$ (with $R \gg 0$ ), where the integral on the RHS of the last equality converges absolutely. However, we shall presently show that for suitable choices of $f$, this integral admits meromorphic continuation to $\mathbb{C}^{r-1}$.

Let us write the torus $T_{E}$ as $T_{0} \times T_{1}$, where $T_{0} \cong F^{\times}$via the coweight $\mu_{0}^{\prime}$, and $T_{1} \cong E^{\times}$ via the coweights $\left\{\mu_{1}^{\prime}, \mu_{2}^{\prime}, \mu_{3}^{\prime}\right\}$. Correspondingly, we shall write the torus element $t$ as $t_{0} t_{1}$. If we conjugate $t_{1}$ to the left across $u_{\alpha} u_{-\beta}$, and change variables, we see that
$Y^{\prime}(\varphi, f, \underline{s})=\int_{T_{0}} \int_{T_{1}} \int_{U_{\alpha} \times U_{-\beta}} \mu_{\underline{s}}\left(t_{0} t_{1}\right) \cdot W_{V_{E}^{\prime}, \psi}\left(t_{1} u_{\alpha} u_{-\beta} t_{0} \cdot \varphi\right) \cdot \overline{L_{\chi_{0}}\left(t_{1} u_{\alpha} u_{-\beta} t_{0} \cdot f\right)} d u_{\alpha} d u_{-\beta} d t_{0} d t_{1}$,
where $\mu_{\underline{s}}$ is the resulting character of $T_{E}$ after these manipulations.
At this point, we need the following lemma which gives an alternative construction of the functional $L_{\chi_{0}}$ on the minimal representation $\Pi_{E}$.

Lemma 9.5.1. Let $P^{\prime \prime}=M^{\prime \prime} N^{\prime \prime}$ be the Heisenberg parabolic subgroup of $G_{E}$ so that the center of $N^{\prime \prime}$ is $Z^{\prime \prime}=U_{\alpha_{0}^{\prime}}=U_{3 \alpha+\beta}$. Let $P_{s s}^{\prime \prime}=M_{s s}^{\prime \prime} N^{\prime \prime}$ be the derived group of $P^{\prime \prime}$. The unique irreducible representation of the Heisenberg group $N^{\prime \prime}$ with central character $\psi$ can be realized on $\mathcal{S}\left(U_{\alpha} \times U_{-\beta}\right)=\mathcal{S}(E \times F)$ and this representation extends uniquely to the Weil representation $\omega_{\psi}$ of $P_{s s}^{\prime \prime}$.
(i) There is a $P^{\prime \prime}$-equivariant injective map

$$
\iota: \Pi_{E} \longrightarrow \operatorname{Ind} d_{P_{s,}^{\prime \prime}}^{P^{\prime \prime}} \omega_{\psi}
$$

(ii) We may realize the latter induced representation on the space of smooth functions on $T_{0} \cong F^{\times}$taking values in $\mathcal{S}\left(U_{\alpha} \times U_{-\beta}\right)$. Thus a function in this space may be denoted by $\phi\left(t_{0} ; u_{\alpha}, u_{-\beta}\right)$. Moreover, the image of $\iota$ contains $\mathcal{S}\left(T_{0}\right) \otimes \mathcal{S}\left(U_{\alpha} \times U_{-\beta}\right)$.
(iii) The linear functional $L_{\chi_{0}}$ is given by

$$
L_{\chi_{0}}(f)=\iota(f)(1 ; 0,0)
$$

for $f \in \Pi_{E}$. Thus, we have:

$$
L_{\chi 0}\left(t_{1} u_{\alpha} u_{-\beta} t_{0} \cdot f\right)=\chi\left(t_{1}\right) \cdot \iota(f)\left(t_{0} ; u_{\alpha}, u_{-\beta}\right),
$$

where $\chi$ is a character of $T_{1}$ which we will not make explicit here.

Proof. We shall give a sketch of the proof. The statements (i) and (ii) follow from the very construction of the minimal representation (cf. [K] and [GS]). To deduce (iii), we observe that via Frobenius reciprocity, $\iota$ gives a $P_{s s}^{\prime \prime}$-equivariant map (the local Fourier-Jacobi map)

$$
F J_{\psi}: \Pi_{E} \longrightarrow \mathcal{S}\left(U_{\alpha} \times U_{-\beta}\right) .
$$

This is explicitly described by

$$
F J_{\psi}(f)\left(u_{\alpha}, u_{-\beta}\right)=\iota(f)\left(1 ; u_{\alpha}, u_{-\beta}\right)
$$

In the $p$-adic case, this is in fact the unique such map, since $\left(\Pi_{E}\right)_{Z^{\prime \prime}, \psi} \cong \omega_{\psi}$. If we let

$$
R=\left\langle Z^{\prime \prime}, U_{3 \alpha+2 \beta}, U_{2 \alpha+\beta}, U_{\alpha+\beta}\right\rangle \subset N_{E},
$$

then the composition of $F J_{\psi}$ with evaluation at $(0,0)$ gives a linear functional $l$ on $\Pi_{E}$ satisfying

$$
l(r \cdot f)=\chi_{0}(r) \cdot l(f) .
$$

In the $p$-adic case, we have $\left(\omega_{\psi}\right)_{R, \chi_{0}} \cong \mathbb{C}$, so that

$$
\left(\Pi_{E}\right)_{R, \chi_{0}} \cong \mathbb{C} .
$$

Thus, $l=L_{\chi_{0}}$ up to scaling. The archimedean case can be finessed from this by a global argument using weak approximation; we omit the details.

In view of the lemma, we have

$$
\begin{equation*}
Y^{\prime}(\varphi, f, \underline{s})=\int_{T_{0}} \int_{U_{\alpha} \times U_{-\beta}} \overline{\iota(f)\left(t_{0} ; u_{\alpha}, u_{-\beta}\right)} \cdot \mu_{\underline{s}}\left(t_{0}\right) \cdot Y^{\prime \prime}\left(u_{\alpha} u_{-\beta} t_{0} \cdot \varphi, \underline{s}\right) d u_{\alpha} d u_{-\beta} d t_{0}, \tag{9.5.2}
\end{equation*}
$$

where

$$
Y^{\prime \prime}(\varphi, \underline{s})=\int_{T_{1}} \mu_{\underline{s}}\left(t_{1}\right) \cdot \chi\left(t_{1}\right) \cdot W_{V_{E}^{\prime}, \psi}\left(t_{1} \cdot \varphi\right) d t_{1} .
$$

Now as we noted in the proof of Prop. 9.1.3(ii) (after the proof of Lemma 9.4.1), an integral of the type defining $Y^{\prime \prime}(\varphi, \underline{s})$ (i.e. the Mellin transform of a Whittaker function) admits meromorphic continuation to $\mathbb{C}^{r-1}(\mathrm{cf} .[\mathrm{JS}, \S 3])$. Further, at a point $s_{0}$ of holomorphy, one can show that the linear functional $\varphi \mapsto Y^{\prime \prime}\left(\varphi, \underline{s_{0}}\right)$ is continuous. Thus the map $\left(t_{0}, u_{\alpha}, u_{-\beta}\right) \mapsto$ $Y^{\prime \prime}\left(t_{0} u_{\alpha} u_{-\beta} \cdot \varphi, \underline{s}\right)$ defines a smooth function.

Now by Lemma 9.5.1(ii), we may take $f$ so that

$$
\overline{\iota(f)\left(t_{0} ; u_{\alpha}, u_{-\beta}\right)}=f_{0}\left(t_{0}\right) \cdot f_{1}\left(u_{\alpha}, u_{-\beta}\right)
$$

for arbitrary $f_{0} \in C_{c}^{\infty}\left(T_{0}\right)$ and $f_{1} \in C_{c}^{\infty}\left(U_{\alpha} \times U_{-\beta}\right)$. Thus, the integrals over $T_{0}$ and $U_{\alpha} \times U_{-\beta}$ converges absolutely and the equation (9.5.2) gives the meromorphic continuation of $Y^{\prime}$.

By the above discussion, we see that if $Y^{\prime}(\varphi, f, \underline{s})$ vanishes for all choices of $f$, then $Y^{\prime \prime}(\varphi, \underline{s})$ vanishes for all choices of $\varphi$. We are thus reduced to showing that there exists $\varphi$ such that

$$
\int_{T_{1}} W_{V_{E}^{\prime}, \psi}\left(t_{1} \cdot \varphi\right) \cdot \mu_{\underline{s}}\left(t_{1}\right) \cdot \chi\left(t_{1}\right) d t_{1}
$$

is absolutely convergent for all $\underline{s}$ and non-zero for a particular $\underline{s}$. For this, we take $\phi \in$ $\mathcal{S}\left(U_{-\alpha}\right)=\mathcal{S}(E)$ and replace $\varphi$ by

$$
\phi * \varphi=\int_{U_{-\alpha}} \phi(u) \cdot \pi(u) \varphi d u .
$$

Then

$$
W_{V_{E}^{\prime}, \psi}\left(t_{1} \cdot(\phi * \varphi)\right)=\widehat{\phi}\left(t_{1}\right) \cdot W_{V_{E}^{\prime}, \psi}\left(t_{1} \cdot \varphi\right)
$$

where $\widehat{\phi}$ is the Fourier transform of $\phi$. We may choose $\phi$ so that $\widehat{\phi}$ is an arbitrary compactly supported function. Then, for a suitable choice of such a $\phi$, the integral over $T_{1}$ converges absolutely for all $\underline{s}$ and can be arranged to be non-zero for any given $\underline{s_{0}}$.

This completes the proof of Prop. 9.1.3(iii).

## 10. Polynomial Invariants

In the case of split groups, Ginzburg has observed that all $L$-functions $L(s, \pi, \rho)$ which are known to be representable by some Rankin-Selberg integrals have the property that the representation $(\rho, V)$ of ${ }^{L} G$ has the property that $\mathbb{C}[V]^{L} G$ is a polynomial ring. This observation appears, for example, in [GH2] (where the word "split" was mistakenly omitted). Ginzburg has confirmed in private communication that his observation does not address the non-split cases.

In any case, the referee of this paper suggested that we investigate if Ginzburg's observation remains valid in our case. Thus, we shall describe briefly the algebras of ${ }^{L} G^{0}$ - and ${ }^{L} G$ invariant polynomials on

$$
V_{\text {Spin }}=V_{\mu_{1}} \oplus V_{\mu_{2}} \oplus V_{\mu_{3}} .
$$

For each $i$, there is a (unique up to scaling) nondegenerate quadratic form $Q_{i}$ on $V_{\mu_{i}}$ which is fixed by ${ }^{L} G^{0}$. We may normalize these so that the action of $S_{3}$ permutes them. It is well-known that

$$
\mathbb{C}\left[V_{\mu_{i}}\right]^{L} G^{0}=\mathbb{C}\left[Q_{i}\right] .
$$

In addition, there is an ${ }^{L} G^{0}$-invariant trilinear linear form $R$ (unique up to scaling) on $V_{\mu_{1}} \times$ $V_{\mu_{2}} \times V_{\mu_{3}}$. In fact, one can show that $\mathbb{C}\left[V_{\text {Spin }}\right]^{L} G^{0}$ is the polynomial algebra generated by these four elements, as may be deduced from the local computations when $E=F \times F \times F$.

In fact, using the geometric description of $V_{\text {Spin }}$ in (2.8), one can write down the above invariants easily. If $\left(x_{1}, x_{2}, x_{3}\right) \in V_{\text {Spin }}=\mathbb{O}^{3}$, then for $i=1,2,3$, we have:

$$
Q_{i}\left(x_{1}, x_{2}, x_{3}\right)=N\left(x_{i}\right) \quad \text { and } \quad R\left(x_{1}, x_{2}, x_{3}\right)=\operatorname{Tr}\left(x_{1} x_{2} x_{3}\right) .
$$

From this description, it is easy to verify that the $Q_{i}$ 's are permuted by $S_{3}$ whereas $R$ is fixed by $S_{3}$. For example, when $\sigma$ is the transposition such that

$$
\sigma:\left(x_{1}, x_{2}, x_{3}\right) \mapsto\left(\overline{x_{2}}, \overline{x_{1}}, \overline{x_{3}}\right),
$$

then

$$
R\left(\sigma\left(x_{1}, x_{2}, x_{3}\right)\right)=\operatorname{Tr}\left(\overline{x_{2} x_{1} x_{3}}\right)=\operatorname{Tr}\left(x_{3} x_{1} x_{2}\right)=R\left(x_{1}, x_{2}, x_{3}\right) .
$$

So for example,

$$
\mathbb{C}\left[V_{\text {Spin }}\right]^{\operatorname{Spin}_{8} \rtimes S_{3}}=\mathbb{C}\left[\Sigma_{1}, \Sigma_{2}, \Sigma_{3}, R\right]
$$

where the $\Sigma_{i}$ 's are the elementary symmetric functions in the $Q_{i}$ 's.
There are thus four possibilities for $\mathbb{C}\left[V_{\text {Spin }}\right]^{L} G$ corresponding to the four possibilities for the image of the Galois group in $S_{3}$ :

- when $E=F \times F \times F$,

$$
\mathbb{C}\left[V_{\text {Spin }}\right]^{L} G=\mathbb{C}\left[Q_{1}, Q_{2}, Q_{3}, R\right]
$$

- when $E=F \times K$,

$$
\mathbb{C}\left[V_{\text {Spin }}\right]^{L} G=\mathbb{C}\left[Q_{1}, Q_{2}+Q_{3}, Q_{2} Q_{3}, R\right] .
$$

- when $E$ is Galois cubic,

$$
\mathbb{C}\left[V_{\text {Spin }}\right]^{L_{G}}=\mathbb{C}\left[Q_{1}, Q_{2}, Q_{3}\right]^{A_{3}}[R],
$$

where $\mathbb{C}\left[Q_{i}\right]^{A_{3}}$ is the direct sum of the symmetric and skew-symmetric polynomials in the $Q_{i}$ 's.

- when $E$ is non-Galois cubic,

$$
\mathbb{C}\left[V_{\text {Spin }}\right]^{L_{G}}=\mathbb{C}\left[\Sigma_{1}, \Sigma_{2}, \Sigma_{3}, R\right] .
$$

Observe that $\mathbb{C}\left[V_{\text {Spin }}\right]^{L} G$ is a polynomial ring except when $E$ is Galois cubic. This suggests that, if Ginzburg's observation were to extend to the non-split case, it is the ${ }^{L} G^{0}$-invariants, rather than the ${ }^{L} G$-invariants, that should play a role.

## 11. Appendix: Tensor Products and Universal Characters

In this section, all groups are over $\mathbb{C}$, so the " $\mathbb{C}$ "'s are suppressed throughout. Thus $G L_{n}$ means $G L_{n}(\mathbb{C})$, etc. For $G$ a connected reductive Lie group, let $R(G)$ denote the representation ring of $G$. Our primary goal here is to give an exposition of some of the results in $[\mathrm{B}-\mathrm{K}-\mathrm{W}]$ from a point of view akin to that of $[\mathrm{K}-\mathrm{T}]$.

It is perhaps worthwhile to comment on a potential source of confusion: the relationship between highest weights and partitions (i.e., nonincreasing sequences of integers of finite length). For classical groups there is a natural identification between a subset of the set of partitions (those with a number of parts bounded suitably in terms of the group) and a subset of the set of weights (for example, for $S p i i_{2 n+1}$, the weights which factor through the projection to $S O_{2 n+1}$ ). Greek letters, $\lambda, \mu, \nu$, etc., are hence used for both partitions and weights. One must therefore be a bit careful about what " $\lambda$ " is in any given sentence: is it a partition, which may have too many parts, and hence not be a weight, or is it a weight, which may not lie in the subset parametrized by partitions?

For $G=G L_{n}$, we consider the subring $R_{+}(G)$ of "polynomial representations," i.e., representations whose characters are polynomials in the coordinates of an element of the torus. (The character of a general representation being the product of one of these times a power of the inverse of the determinant.) The highest weight of an element of $R_{+}(G)$ is a partition with at most $n$ parts, $\lambda$, the value at $\operatorname{diag}\left(t_{1}, \ldots, t_{n}\right)$ being $\prod_{i=1}^{n} t_{i}^{\lambda_{i}}$. Here, if $\lambda$ has fewer than $n$ parts, we add terminal zeros.

Since the map that sends a representation to its character is injective, we shall generally speak in terms of characters, rather than representations. We introduce a little more general notation. For $G$ as above, and $\lambda$ a dominant weight of $G$ we let $\Gamma_{G}(\lambda)$ denote the irreducible finite dimensional representation of $G$ with highest weight $\lambda$, and $\chi_{G}(\lambda)$ its character.
11.1. Review and Summary. We review the method of universal characters as discussed in $[\mathrm{K}-\mathrm{T}]$. For each $n$, let $\Lambda_{n}=R_{+}\left(G L_{n}\right)=\mathbb{Z}\left[t_{1}, \ldots, t_{n}\right]^{\mathfrak{G}_{n}}$. For $m>n$, let $\tilde{\rho}_{m, n}: \Lambda_{m} \rightarrow \Lambda_{n}$ be the ring homomorphism that sends $t_{i}$ to $t_{i}$ if $i \leq n$ and 0 if $i>n$. Then $\left(\Lambda_{n}, \tilde{\rho}_{m, n}\right)$ is a projective system. The projective limit (in the category of $\mathbb{Z}_{\geq 0}$-graded algebras), $\Lambda$, with maps $\tilde{\rho}_{n}$ to each $\Lambda_{n}$, was defined by MacDonald [M] and is called the universal character ring.

For each $n$, we define the elementary symmetric polynomials, $e_{k}^{n}(t)$ by

$$
\prod_{i=1}^{n}\left(1+t_{i} x\right)=\sum_{k=-\infty}^{\infty} e_{k}^{n}(t) x^{k} .
$$

In particular, $e_{k}^{n}(t)=0$ if $k<0$ or $k>n$. Then $\mathbb{Z}\left[t_{1}, \ldots, t_{n}\right]^{\mathfrak{G}_{n}}=\mathbb{Z}\left[e_{1}^{n}, \ldots, e_{n}^{n}\right]$. Clearly $\tilde{\rho}_{m, n}\left(e_{k}^{m}\right)=e_{k}^{n}$, so we get elements $e_{k}$ of $\Lambda$. Similarly, we consider $p_{i}^{n}$ defined by

$$
\prod_{i=1}^{n}\left(1-x t_{i}\right)^{-1}=\sum_{i=-\infty}^{\infty} p_{i}^{n}(t) x^{i}
$$

so that $p_{i}^{n}=0$ for $i<0$ but not for $i$ large. Then $\tilde{\rho}_{m, n}\left(p_{i}^{m}\right)=p_{i}^{n}$, and so once again we get elements of $\Lambda$.

We have the classical determinantal expressions

$$
\begin{equation*}
\chi_{G L_{n}}(\lambda)=\left|p_{\lambda_{i}+i-j}^{n}\right|=\left|e_{\lambda_{i}^{\prime}+i-j}^{n}\right| . \tag{11.1.1}
\end{equation*}
$$

Here and throughout $|f(i, j)|$ denotes the determinant of the square matrix whose $i, j$ entry is $f(i, j)$, and $\lambda^{\prime}$ is the transpose of the partition $\lambda$. Note that in each expression, the only dependence of $n$ is the superscripts on the $p_{i}$ 's (resp. $e_{i}$ 's). In particular, the size of the determinant is equal to the number of nonzero parts of $\lambda$ (resp. $\lambda^{\prime}$ ) and does not depend on $n$. Making use of these, together with the elements $p_{i}, e_{i} \in \Lambda$ above, we define $\chi_{G L}(\lambda) \in \Lambda$. For each $n$ we have the natural projection $\pi_{n}: \Lambda \rightarrow \Lambda_{n}$. Then clearly $\pi_{n}\left(\chi_{G L}(\lambda)\right)=\chi_{G L_{n}}(\lambda)$ if $\lambda$ has at most $n$ parts. What is more, if $\lambda$ has more than $n$ parts, then the top row of the second determinantal expression in (11.1.1) is all zeros, so that $\pi_{n}\left(\chi_{G L}(\lambda)\right)=0$.

The universal characters $\chi_{G L}(\lambda)$ are a $\mathbb{Z}$-basis for $\Lambda$, and the corresponding structure constants are the Littlewood-Richardson coefficients $L R_{\mu, \nu}^{\lambda}$. See [F-H, p. 455-56]. Thus,

$$
\chi_{G L}(\mu) \chi_{G L}(\nu)=\sum_{\lambda} L R_{\mu, \nu}^{\lambda} \chi_{G L}(\lambda),
$$

where the sum is over all partitions $\lambda$, whereas

$$
\chi_{G L_{n}}(\mu) \chi_{G L_{n}}(\nu)=\sum_{\lambda} L R_{\mu, \nu}^{\lambda} \chi_{G L_{n}}(\lambda),
$$

where now the sum is only over those partitions with at most $n$ parts.
We turn now to $S p_{2 n}$. Highest weights are once again naturally identified with partitions of length $\leq n$. We have

$$
\chi_{S p_{2 n}}(\lambda)=\left|p_{\lambda_{i}+i-j}^{2 n}+\delta_{j, 1} p_{\lambda_{i}-i-j+2}^{2 n}\right|,
$$

where $\delta_{i, j}$ is a Kronecker $\delta$ and $p_{i}^{2 n}$ has been restricted from the torus of $G L_{2 n}$ to that of $S p_{2 n}$. Once again this is nearly independent of $n$ and allows us to define $\chi_{S p}(\lambda) \in \Lambda$. However, it is no longer true that the determinant simply vanishes when $\lambda$ has more than $n$ parts. The universal characters $\chi_{S p}(\lambda)$ give a second $\mathbb{Z}$-basis for $\Lambda$. The corresponding structure constants may be given in terms of the Littlewood-Richardson coefficients, as in [Ko], Theorem 7.5. This has a nice expression if we adopt one of the notational innovations in $[\mathrm{B}-\mathrm{K}-\mathrm{W}]$, namely the definition

$$
\begin{equation*}
\lambda / \zeta:=\sum_{\substack{\xi \\ 45}} L R_{\zeta, \xi}^{\lambda} \xi \tag{11.1.2}
\end{equation*}
$$

Let $z, x$ and $\ell$ denote the sums of the parts of $\zeta, \xi$ and $\lambda$ respectively. The coefficient $L R_{\zeta, \xi}^{\lambda}$ counts the number of ways to add the $z$ boxes to the Young diagram of $\xi$ (or $x$ to that of $\zeta)$ so as to obtain the Young diagram of $\lambda$, subject to certain conditions. In particular, it is zero for $x+z \neq \ell$, so the above sum is effectively finite, and zero for almost all $\zeta$. With this notation, we have

$$
\begin{equation*}
\chi_{S p}(\lambda) \chi_{S p}(\mu)=\sum_{\zeta} \chi_{S p}((\lambda / \zeta) \cdot(\mu / \zeta)), \tag{11.1.3}
\end{equation*}
$$

where the interior product is given by the Littlewood-Richardson rule. Here we extend the definition of $\chi_{S p}$ to finite sums of partitions by $\mathbb{Z}$-linearity. Hence

$$
\begin{equation*}
\chi_{S p_{2 n}}(\lambda) \chi_{S p_{2 n}}(\mu)=\pi_{S p_{2 n}}\left(\sum_{\zeta} \chi_{S p}((\lambda / \zeta) \cdot(\mu / \zeta))\right) . \tag{11.1.4}
\end{equation*}
$$

Here $\pi_{S p_{2 n}}$ is the projection from $\Lambda$ to $R\left(S p_{2 n}\right)$ obtained by projecting to $\Lambda_{2 n} \subset R\left(G L_{2 n}\right)$ and then restricting to $S p_{2 n}$. This formula essentially goes back to Littlewood and Newell. Only the interpretation in terms of universal characters is new.

As noted above, it is not the case that $\pi_{S p_{2 n}}$ gives the naive projection onto the set of partitions with at most $n$ parts. A nice description of the projection and its kernel is given in [Ko, Section 9]. For purposes of our application to $S p_{6}$, we only need to know that

$$
\pi_{S p_{6}}\left(\chi_{S p}(\lambda)\right)=\left\{\begin{array}{l}
\chi_{S p_{6}}(\lambda) \text { if } \lambda \text { has at most } 3 \text { parts; } \\
0, \text { if } \lambda \text { has exactly } 4 \text { parts. }
\end{array}\right.
$$

11.2. New Material. We turn now to the formulae in section 7 of $[B-K-W]$. These are derived by similar means, but in contrast to the above, they relate characters of representations of classical groups to those of general linear groups not by restricting, say, from $G L_{2 n}$ to $S O_{2 n}$, but rather by restricting from $S O_{2 n}$ to $G L_{n}$ embedded as the Levi factor of the Siegel parabolic. The symplectic and odd orthogonal cases are also considered in section 7 of [B-K-W], but it is only the even orthogonal case where the method there really offers substantial improvement over the method of section 5 of that paper and the previous section of these notes.

It is convenient to think of $S O_{2 n}$ as a subset of $G L_{2 n}$, with the maximal torus of diagonal elements $t=\operatorname{diag}\left(t_{1}, \ldots, t_{n}, t_{n}^{-1}, \ldots, t_{1}\right)$, and identify the partition $\lambda$ with the map $t \mapsto \prod_{i} t_{i}^{\lambda_{i}}$. Indeed, we have already done so in referring to "the Siegel parabolic." Of course, there is another identification which works just as well (using the other parabolic) and in the case $n=4$, which is our primary interest here, there are still more possibilities arising from the additional symmetry of the Dynkin diagram. We fix once and for all one of these identifications. When $n \neq 4$, the highest weight of the unique projection to $\mathrm{SO}_{2 n}$ is identified with the partition 1 having one part. When $n=4$, we will also refer to the representation whose highest weight has been identified with this partition as "the" projection. Because of our choice of identification, it is no longer on equal footing with the other two.

There are three matters which must be addressed.
(i) The first is that the restriction of a representation of $S O_{2 n}$ to $G L_{n}$ will not be contained in $\Lambda_{n}=R_{+}\left(G L_{n}\right)$. Hence we can anticipate the need for some larger "universal" ring which covers all of $R\left(G L_{n}\right)$.
(ii) The second is that we are not satisfied with $R\left(S O_{2 n}\right)$, but want all of $R\left(S p i n_{2 n}\right)$. It will turn out that for the cases we need here this second issue may to a certain extent be ducked. We must replace $G L_{n}$ by its inverse image in $\operatorname{Spin}_{2 n}$, which is a double cover $\widetilde{G L}_{n}$, but as we shall see later, the double cover enters the computation only in a trivial way. Alternatively, everything could be phrased in terms of the Lie algebras.
(iii) The third issue is that, while the product of two elements of $R\left(\widetilde{G L}_{n}\right)$, each of which is a restriction from $R\left(\operatorname{Spin}_{2 n}\right)$ is certainly the restriction of something, we need a way to recognize what element of $R\left(\operatorname{Spin}_{2 n}\right)$ it is a restriction of. This is handled by the following Lemma, akin to Brauer's method, which is very pretty in its own right. (It will be applied with $G=\operatorname{Spin}_{2 n}$ and $H=\widetilde{G L}_{n}$ the Levi subgroup of the Siegel parabolic.) It is due to King [Ki].

To state the lemma we need some notation, as before, let $G$ be a connected reductive Lie group. Fix maximal torus $T$ and a choice of positive roots, and let $H$ be a standard Levi (there is a more general formulation, but restricting to standard Levi subgroups is sufficient for our purposes here and simplifies certain things a bit). Then we obtain a set of positive roots for $H$ with respect to $T$, and the notions of $G$-dominance and $H$-dominance for elements of the lattice $L$ of weights. Let $L_{G}^{+}$denote the set of $G$-dominant weights and $L_{H}^{+}$the set of $H$-dominant ones. Then $L_{G}^{+} \subset L_{H}^{+}$since $G$-dominance is more restrictive than $H$-dominance.

For $\lambda \in L$ and $t \in T$, let

$$
A_{\lambda}^{G}(t)=\sum_{w \in W_{G}} \operatorname{sgn}(w) t^{w \lambda}
$$

and

$$
\chi_{G}(\lambda)=A_{\lambda+\rho_{G}} / A_{\rho_{G}} .
$$

For $\lambda G$-dominant, this agrees with the previous definition as the character of $\Gamma_{G}(\lambda)$. For $\lambda$ not $G$-dominant, it is equal to $\operatorname{sgn}(w) \chi_{\tau}$ if $\tau=w\left(\lambda+\rho_{G}\right)-\rho_{G}$ is $G$-dominant, and is zero otherwise.

Lemma 11.2.1. With this notation, we have

$$
\chi_{G}(\lambda) \chi_{G}(\mu)=\sum_{\tau \in L_{H}^{+}}\left(\sum_{\eta} n_{\eta}^{\lambda} c(H)_{\eta, \mu}^{\tau}\right) \chi_{G}(\tau),
$$

where $n_{\eta}^{\lambda}$ is the multiplicity of $\Gamma_{H}(\eta)$ in the restriction $\Gamma_{G}(\lambda)$ to $H$, and $c(H)_{\lambda, \mu}^{\tau}$ is the multiplicity of $\Gamma_{H}(\tau)$ in $\Gamma_{H}(\lambda) \otimes \Gamma_{H}(\mu)$.

Note that this decomposition is in terms of $\chi_{G}$ 's, but the $\tau$ 's are only $H$-dominant. Before we proceed to the proof, let us note that

$$
\chi_{G}(\lambda) \chi_{H}(\mu)=\sum_{\tau \in L_{H}^{+}}\left(\sum_{\eta} n_{\eta}^{\lambda} c(H)_{\eta, \mu}^{\tau}\right) \chi_{H}(\tau),
$$

as is immediate from the definitions. So the lemma may be interpreted as follows. Suppose we have a formula for the restriction from $G$ to $H$ and a method of computing tensor products of representations of $H$. Then we may compute the product of two irreducible representations
of $G$ as follows: restrict one of them to $H$, take the product of this restriction with the irreducible representation of $H$ that has the same highest weight as the other, and then simply replace each of the $\Gamma_{H}$ 's in the answer by the corresponding $\Gamma_{G}$. This is precisely the method employed in [B-K-W] section 7, the requisite branching rules having been obtained in earlier work.

We now remark why it is possible to essentially duck the issue of the difference between $R\left(S O_{2 n}\right)$ and $R\left(\operatorname{Spin}_{2 n}\right)$ in the cases we are dealing with. This is because we never need to compute a product where both of the characters involved are in $R\left(\operatorname{Spin}_{2 n}\right)-R\left(S O_{2 n}\right)$. When one of them is, we make the corresponding weight " $\mu$ " rather than " $\lambda$." Suppose $\lambda$ is a weight of $S p i i_{2 n}$, but not of $S O_{2 n}$. Then either it is of the form $\left(\frac{1}{2}, \ldots, \frac{1}{2}, \frac{1}{2}\right)+p$ with $p$ a partition, or it is associated to a weight of this form by the automorphism that reverses the last two fundamental weights. It is enough to consider the case when $\lambda=\left(\frac{1}{2}, \ldots, \frac{1}{2}, \frac{1}{2}\right)+p$. It is not a weight of $G L_{n}$, but of the double cover $\widetilde{G L}_{n}$. But since

$$
\chi_{\widetilde{G L}_{n}}\left(\left(\frac{1}{2}, \ldots, \frac{1}{2}, \frac{1}{2}\right)+p\right)=\operatorname{det}^{\frac{1}{2}} \chi_{G L_{n}}(p)
$$

the double cover comes in only in a trivial way.
Proof of Lemma 11.2.1: Let $\rho_{G}$ denote half the sum of the positive roots of $G$ and $\rho_{H}$ half the sum of the positive roots in $H$. Let $W_{G}$ denote the Weyl group of $G$ with respect to $T$. Note that the Weyl group $W_{H}$ of $H$ is naturally identified with a subgroup of the Weyl group $W_{G}$ of $G$, and that $\rho_{G}-\rho_{H}$ is $W_{H}$-stable.

Let $m_{\lambda}^{G}(\nu)$ denote the multiplicity of the weight $\nu$ in $\Gamma_{G}(\lambda)$, and define $m_{\lambda}^{H}(\nu)$ similarly. Then

$$
\begin{equation*}
\chi_{G}(\lambda) \chi_{G}(\mu)=\sum_{\nu} m_{\lambda}^{G}(\nu) \chi_{G}(\nu+\mu)=\sum_{\eta} \sum_{\nu} n_{\eta}^{\lambda} m_{\eta}^{H}(\nu) \chi_{G}(\nu+\mu) . \tag{11.2.2}
\end{equation*}
$$

For $\nu$ a weight we define $|\nu|_{G}$ and $\frac{\nu}{|\nu|_{G}}$ as follows. If the stabilizer of $\nu+\rho_{G}$ in $W_{G}$ is trivial, then $|\nu|_{G}$ is the unique dominant weight $\tau$ such that $\tau+\rho_{G}=w\left(\nu+\rho_{G}\right)$ for some $w \in W_{G}$ and $\frac{\nu}{|\nu|_{G}}$ is the sign of this $w$. If the stabilizer of $\nu+\rho_{G}$ in $W_{G}$ is not trivial, then $\frac{\nu}{|\nu|_{G}}=0$, and it does not matter what $|\nu|_{G}$ is (since it is multiplied by zero).

Then from the first equality above, we obtain

$$
c(G)_{\lambda, \mu}^{\tau}=\sum_{|\nu+\mu|_{G}=\tau} m_{\lambda}^{G}(\nu) \frac{\nu+\mu}{|\nu+\mu|_{G}}=\sum_{w \in W_{G}} \operatorname{sgn}(w) m_{\lambda}^{G}\left(w\left(\tau+\rho_{G}\right)-\mu-\rho_{G}\right)
$$

Similarly,

$$
c(H)_{\mu, \eta}^{\tau}=\left(\sum_{\nu:|\mu+\nu|_{H}=\tau} m_{\eta}^{H}(\nu) \frac{\mu+\nu}{|\mu+\nu|_{H}}\right) .
$$

Now $|\mu+\nu|_{H}=\tau$ means that $w\left(\mu+\nu+\rho_{H}\right)=\tau+\rho_{H}$ for some $w \in W_{H}$ in which case $w\left(\mu+\nu+\rho_{G}\right)=\tau+\rho_{G}$, since $\rho_{G}-\rho_{H}$ is $W_{H}$-stable. Thus (11.2.2) equals

$$
\sum_{\tau \in L_{H}^{+}} \sum_{\eta \in L} n_{\eta}^{\lambda} \sum_{\nu:|\mu+\nu|_{H}=\tau} m_{\eta}^{H}(\nu) \frac{\mu+\nu}{|\mu+\nu|_{H}} \chi_{G}(\tau),
$$

and the result follows.

Remark: One can, and King does, formulate and prove a more general result applicable to the case $G=G_{2}, H=S L_{3}$. In this case $\rho_{G}-\rho_{H}$ is no longer $W_{H}$-stable, so the expression is more complicated:

$$
\chi_{G}(\lambda) \chi_{G}(\mu)=\sum_{\tau \in L_{H}^{+}} \sum_{\eta \in L} n_{\eta}^{\lambda} c(H)_{\mu+\rho_{G}-\rho_{H}, \eta}^{\tau} \chi_{G}\left(\tau-\rho_{G}+\rho_{H}\right) .
$$

11.3. An analog of $\Lambda$. As noted above, in order to place the branching rule method of King and his collaborators within a framework analogous to that of $[\mathrm{K}-\mathrm{T}]$, we need an analog of $\Lambda$ that surjects onto all of $R\left(G L_{n}\right)$ for each $n$, not only $R_{+}\left(G L_{n}\right)$.

Let $R_{n}=\mathbb{Z}\left[t_{1}, \ldots, t_{n}, u_{1}, \ldots, u_{n}\right]^{\mathfrak{S}_{n} \times \mathfrak{S}_{n}}$, where $\mathfrak{S}_{n}$ is the symmetric group, the first acting by permuting the $t$-indices and the second permutes the $u$-indices. Let $\pi_{n}$ be the map $R_{n} \rightarrow$ $R\left(G L_{n}\right)$, defined by $u_{i} \mapsto t_{i}^{-1}$. Also, for $m>n$ there is a surjective map $\rho_{m, n}: R_{m} \rightarrow R_{n}$ defined by $t_{i} \mapsto t_{i}, u_{i} \mapsto u_{i}, i \leq n$ and $t_{i} \mapsto 0, u_{i} \mapsto 0, i>n$. Clearly, $\pi_{n} \circ \rho_{m, n}$ does not factor through $\pi_{m}$. However, $\left(R_{n}, \rho_{m, n}\right)$ forms a projective system. We let $R$ denote the projective limit in the category of $\mathbb{Z}_{\geq 0} \times \mathbb{Z}_{\geq 0}$ graded algebras, and let $\rho_{n}$ denote the natural map $R \rightarrow R_{n}$.

We denote the $n$-tuples $\left(t_{1}, \ldots, t_{n}\right),\left(u_{1}, \ldots, u_{n}\right),\left(t_{1}^{-1}, \ldots, t_{n}^{-1}\right)$ etc., more briefly by $t, u, t^{-1}$, etc.

Lemma 11.3.1. For each $n$, we have $R_{n}=\mathbb{Z}\left[e_{1}^{n}(t), \ldots, e_{n}^{n}(t), e_{1}^{n}(u), \ldots, e_{n}^{n}(u)\right]$.
From this it is clear that the map $\pi_{n}$ is surjective, since $R\left(G L_{n}\right)$ is generated by $e_{i}^{n}(t), i=$ $1, \ldots, n$ and $\operatorname{det}^{-1}=\pi_{n}\left(e_{n}^{n}(u)\right)$.

Proof: This is clear: $\mathbb{Z}\left[t_{1}, \ldots, t_{n}\right]^{\mathfrak{S}_{n}}=\mathbb{Z}\left[e_{1}^{n}(t), \ldots, e_{n}^{n}(t)\right]$ (likewise with $u$ 's) and

$$
R_{n}=\mathbb{Z}\left[t_{1}, \ldots, t_{n}\right]^{\mathfrak{S}_{n}} \otimes_{\mathbb{Z}} \mathbb{Z}\left[u_{1}, \ldots, u_{n}\right]^{\mathfrak{S}_{n}} .
$$

Since $\rho_{m, n}\left(e_{k}^{m}(t)\right)=e_{k}^{n}(t)$ (likewise with $u$ ), we obtain elements of $R$, which we denote by $e_{k,+}$ and $e_{k,-}$, such that $\rho_{n}\left(e_{k,+}\right)=e_{k}^{n}(t)$, and $\rho_{n}\left(e_{k,-}\right)=e_{k}^{n}(u)$ for each $n$. The notation $e_{k,-}$ is, of course, motivated by the fact that $\pi_{n} \circ \rho_{n}\left(e_{k,-}\right)=e_{k}^{n}\left(t^{-1}\right)$.

It follows from Lemma 11.3.1 that $R=\mathbb{Z}\left[\left\{e_{k, \pm}\right\}\right]$.
Lemma 11.3.2. For every element $\chi$ of $R$, there is an $N$ such that if $n>N$, then $\chi$ is not in the kernel of $\pi_{n} \circ \rho_{n}$.

Proof: Clearly, for each $n$, the map $\rho_{n}$ is injective on $\mathbb{Z}\left[\left\{e_{k, \pm}: k \leq n\right\}\right]$. The kernel of $\pi_{n}$ is generated by $\left\{e_{n}^{n}(t) e_{i}^{n}(u)-e_{n-i}^{n}(t): i=0, \ldots, n\right\}$. Given $\chi \in R$, we define $K_{1}$ to be the largest $k$ such that $e_{k,+}$ appears in the expression of $\chi$ in terms of the $e_{i, \pm}$, and $K_{2}$ to be the largest $k$ such that $e_{k,-}$ appears. Then $N=K_{1}+K_{2}+1$ is sufficient.

Next, we wish to define an element of $R$ associated to any pair of partitions $\lambda$ and $\mu$. This is given in terms of the transpose partitions $\lambda^{\prime}$ and $\mu^{\prime}$, and it will be convenient to denote $\lambda_{1}$, which is equal to the number of parts in $\lambda^{\prime}$, more briefly by $r$. Similarly, we let $\mu_{1}=s$.

Then, let

$$
\chi_{G L}(\bar{\mu} ; \lambda)=\left|\begin{array}{ccc}
e_{\mu_{s}^{\prime},-} & \ldots & e_{\mu_{s}^{\prime}-r-s+1,-} \\
\vdots & \ddots & \vdots \\
e_{\mu_{1}^{\prime}+s-1,-} & \cdots & e_{\mu_{1}^{\prime}-r,-} \\
e_{\lambda_{1}^{\prime}-s,+} & \ldots & e_{\lambda_{1}^{\prime}+r-1,+} \\
\vdots & \ddots & \vdots \\
e_{\lambda_{r}^{\prime}-r-s+1,+} & \cdots & e_{\lambda_{r}^{\prime},+}^{\prime}
\end{array}\right|
$$

Let $n$ be at least $\lambda_{1}^{\prime}+\mu_{1}^{\prime}$. Then, making use of the identity $\operatorname{det}^{-1}(t) e_{k}^{n}(t)=e_{n-k}\left(t^{-1}\right)$, it is easy to check that

$$
\pi_{n} \circ \rho_{n}\left(\chi_{G L}(\bar{\mu} ; \lambda)\right)=\operatorname{det}^{-s} \chi_{G L_{n}}(\tau(\lambda, \mu)),
$$

where $\tau$ is the partition such that $\tau(\lambda, \mu)^{\prime}=\left(n-\mu_{s}^{\prime}, \ldots, n-\mu_{1}^{\prime}, \lambda_{1}^{\prime}, \ldots, \lambda_{r}^{\prime}\right)$. (This $\tau(\lambda, \mu)$ may be visualized as follows: start with the Young diagram of $\lambda$. Add $\mu_{1}$ columns of length $n$ at the left side. Then take the Young diagram of $\mu$, rotate it 180 degrees, and subtract it from the $n \times \mu_{1}$ rectangle.) So, $\pi_{n} \circ \rho_{n}$ maps $\left\{\chi_{G L}(\bar{\mu} ; \lambda) \mid \lambda_{1}^{\prime}+\mu_{1}^{\prime} \leq n\right\}$ bijectively onto the set of characters of irreducible representations of $G L_{n}$.

Observe that $\Lambda$ is naturally identified with a subring of $R, \Lambda_{n}$ with a subring of $R_{n}$, and the restriction of $\rho_{m, n}$ is $\tilde{\rho}_{m, n}$ for each $m, n$. Also, when $\mu$ is the empty partition, $\chi_{G L}(\bar{\mu} ; \lambda)$ is precisely the universal character $\chi_{G L}(\lambda) \in \Lambda$ defined above.

We may now offer a remark about why we do not simply generate $R\left(G L_{n}\right)$ by $\left\{e_{i}^{n}\right.$ : $i=1, \ldots, n\}$ and $\operatorname{det}^{-1}$. Let $\lambda$ be a partition. The expression for $\chi_{S O_{2 n}}(\lambda)$ as a sum of characters of $G L_{n}$ is stable (in the sense that it is independent of $n$ sufficiently large) when these characters are expressed as $\pi_{n} \circ \rho_{n} \chi_{G L}(\bar{\mu} ; \nu)$. On the other hand, the expression in terms of $\operatorname{det}^{-s} \chi_{G L_{n}}(\tau(\nu, \mu))$ does not stabilize as $n \rightarrow \infty$. Indeed, the partition $\tau(\nu, \mu)$ is different for each value of $n$.
11.4. A review of our application. With the background established, let us go back over the application to the local computations when $E=F \times F \times F$. We must compute the product of two characters. One of them is the character of a representation whose highest weight not only is a partition, but is one with strictly less than $n$ parts. Let us call this partition $\lambda$.

We first need a formula for the restriction of $\chi_{S O_{2 n}}(\lambda)$ from $S O_{2 n}$ to $G L_{n}$. Once again the "/" notation of (11.1.2) allows for a nice formulation. The relevant formula, which appears in [B-K-W] as (A4) on p. 1586 is, in our notation:

$$
\begin{equation*}
\chi_{S O_{2 n}}(\lambda)=\sum_{\xi, \beta} \pi_{n} \circ \rho_{n}\left(\chi_{G L_{n}}(\bar{\xi} ; \lambda /(\xi \cdot \beta))\right), \tag{11.4.1}
\end{equation*}
$$

where $\lambda$ has at most $n$ parts, $\xi$ is summed over all partitions, and $\beta$ is summed only over partitions such that each part appears with even multiplicity. This is a finite sum, since there are only finitely many pairs $(\xi, \beta)$ such that $\lambda /(\xi \cdot \beta)$ is nonzero. But note that this holds only for those $\lambda$ which (a) are partitions, as not all weights of $S O_{2 n}$, much less those of $\operatorname{Spin}_{2 n}$ are, and (b) as partitions, have strictly fewer than $n$ parts,

Next, we take the product of (11.4.1) with the character of $\widetilde{G L}_{n}$ which has the same highest weight as our other representation of $\operatorname{Spin}_{2 n}$. Let us assume that this weight is $\mu+(\varepsilon, \ldots, \varepsilon, \varepsilon)$,
where $\mu$ is a partition, and $\varepsilon=0$ or $\frac{1}{2}$. Then this character is $\chi_{G L_{n}}(\mu) \operatorname{det}^{\varepsilon}$. So, we obtain

$$
\begin{equation*}
\chi_{S O_{2 n}}(\lambda) \chi_{\widetilde{G L}_{n}}\left(\mu+\varepsilon^{n}\right)=\operatorname{det}^{\varepsilon} \pi_{n} \circ \rho_{n}\left(\sum_{\xi, \beta} \chi_{G L}(\bar{\xi} ; \lambda /(\xi \cdot \beta)) \chi_{G L}(\mu)\right) . \tag{11.4.2}
\end{equation*}
$$

We now need to evaluate these products of characters of $G L_{n}$, which we wish to do in $R$ and then project. The product in $R$ is described by formula (5.2) p. 1571 of [B-K-W]. In our notation:

$$
\chi_{G L}(\bar{\mu} ; \lambda) \chi_{G L}(\bar{\kappa} ; \nu)=\sum_{\sigma, \tau} \chi_{G L} \overline{(\overline{(\mu / \sigma) \cdot(\kappa / \tau)} ;(\lambda / \tau) \cdot(\nu / \sigma)) .}
$$

Plugging this into (11.4.2) we get

$$
\chi_{S O_{2 n}}(\lambda) \chi_{\widetilde{G L_{n}}}\left(\mu+\varepsilon^{n}\right)=\operatorname{det}^{\varepsilon} \sum_{\eta, \zeta, \beta} \pi_{n} \circ \rho_{n}\left(\chi_{G L}(\bar{\eta} ;(\lambda /(\zeta \eta \beta)) \cdot(\mu / \zeta))\right),
$$

where we have extended $\chi_{G L}$ to pairs of finite formal sums of partitions by $\mathbb{Z}$-linearity. Let us now define, for each $n$ a map $T_{\widetilde{G L}_{n} \rightarrow \operatorname{Spin}_{2 n}}$ from $R\left(\widetilde{G L}_{n}\right)$ to $R\left(\operatorname{Spin}_{2 n}\right)$ which is $\mathbb{Z}$-linear and sends $\chi_{\widetilde{G L_{n}}}(\lambda)$ to $\chi_{S_{p i i_{2 n}}}(\lambda)$ for each weight $\lambda$. Then we get

$$
\begin{align*}
& \chi_{S O_{2 n}}(\lambda) \chi_{S p i n_{2 n}}\left(\mu+\varepsilon^{n}\right)=\sum_{\xi, \beta} T_{\widetilde{G L}_{n} \rightarrow S p i n_{2 n}}\left(\operatorname{det}^{\varepsilon} \chi_{G L_{n}}(\bar{\xi} ; \lambda /(\xi \cdot \beta)) \chi_{G L_{n}}(\mu)\right)  \tag{11.4.3}\\
& \quad=\sum_{\eta, \zeta, \beta} T_{\widetilde{G L}_{n} \rightarrow S p i n_{2 n}}\left(\operatorname{det}^{\varepsilon} \pi_{n} \circ \rho_{n}\left(\chi_{G L}(\bar{\eta} ;(\lambda /(\zeta \eta \beta)) \cdot(\mu / \zeta))\right)\right) \tag{11.4.4}
\end{align*}
$$

All that is needed now to compute $\chi_{S O_{2 n}}(\lambda) \chi_{G L_{n}}(\mu)$ is an efficient algorithm for reducing $\chi_{G L}$ modulo the kernel of $\pi_{n} \circ \rho_{n}$. This is given in the next subsection. The answer is always either 0 or $\pm \chi_{G L_{n}}\left(\bar{\tau}^{\circ} ; \nu^{\circ}\right)$ for a single pair of partitions $\tau^{\circ}, \nu^{\circ}$ with at most $n$ total parts.
11.5. Reduction Modulo $\operatorname{Ker}\left(\pi_{n} \circ \rho_{n}\right)$. In this section we show how, for any pair of partitions $(\bar{\mu} ; \lambda)$ with any numbers of parts, to reduce $\chi_{G L}(\bar{\mu} ; \lambda)$ modulo the kernel of $\pi_{n} \circ \rho_{n}$, obtaining either zero, or a pair ( $\bar{\mu}^{\circ} ; \lambda^{\circ}$ ) with at most $n$ total parts, such that $\pi_{n} \circ \rho_{n}\left(\chi_{G L}(\bar{\mu} ; \lambda)\right)= \pm \chi_{G L_{n}}\left(\bar{\mu}^{\circ} ; \lambda^{\circ}\right)$, with an explicit description of the sign. This corresponds to the $U_{n}, S U_{n}$ modification rule in Table 3 of [B-K-W], although it may be necessary to apply that rule more than once to obtain the answer we describe below. By contrast, the modification rules in Tables 4 and 5 correspond to finding the dominant weight $|\nu|_{G}$ associated to $\nu$ by the action of the Weyl group, shifted by $\rho_{G}$.

For any infinite sequence $a=a_{1}, a_{2}, \ldots, a_{m}, \ldots$ of integers with the property that $a_{m}=$ $1-m$ for $m>N$, we define $E(a)$ to be the $N \times N$ determinant whose $(i, j)$ entry is $e_{a_{i}+j-1}$, and define $E^{n}(a)$ analogously with $e^{n}$ 's. Observe that this definition is independent of the choice of specific $N$ having the requisite property. If $\ell$ is an infinite sequence of integers such that $\ell_{m}=0$ for all $m$ sufficiently large, we follow $[\mathrm{Ko}]$ in defining the associated $\beta$-sequence $\beta(\ell)$ by $\beta(\ell)_{m}=\ell_{m}+1-m$. Observe that if $\lambda$ is a partition, then $\chi_{G L}(\lambda)=E\left(\beta\left(\lambda^{\prime}\right)\right)$, where we have identified the partition $\lambda^{\prime}$ with the infinite sequence obtained by appending zeros at the end.

We now consider $\pi_{n} \circ \rho_{n}\left(\chi_{G L_{n}}(\bar{\mu} ; \lambda)\right)$. Making use of the identity $e_{n}^{i}\left(t^{-1}\right)=e_{n}^{n}(t)^{-1} e_{n-i}^{n}(t)$, we find that

$$
\pi_{n} \circ \rho_{n}\left(\chi_{G L_{n}}(\bar{\mu} ; \lambda)\right)=\left(e_{n}^{n}\right)^{-s} E^{n}(\beta(a(\lambda, \mu)))
$$

where

$$
a(\lambda, \mu)=\left(n-\mu_{s}^{\prime}, \ldots, n-\mu_{1}^{\prime}, \lambda_{1}^{\prime}, \ldots, \lambda_{r}^{\prime}\right)
$$

Here, we have once again denoted $\lambda_{1}$ and $\mu_{1}$ by $r$ and $s$, as is convenient to do when they show up as the number of parts of the transpose partitions $\lambda^{\prime}$ and $\mu^{\prime}$. When $n-\mu_{1}^{\prime} \geq \lambda_{1}^{\prime}$ this expresses $\pi_{n} \circ \rho_{n}\left(\chi_{G L_{n}}(\bar{\mu} ; \lambda)\right.$ as the character of a polynomial representation of $G L_{n}$ times a power of $\operatorname{det}^{-1}$.

Now suppose that $a(\lambda, \mu)$ is not a nonincreasing sequence (i.e., that $\mu_{1}^{\prime}+\lambda_{1}^{\prime}>n$.) Denote $\beta(a(\lambda, \mu))$ more briefly by $\beta$. If there exists $i \neq j$ such that $\beta_{i}=\beta_{j}$, then $E(\beta)=0$. Otherwise, we may rearrange the terms to obtain a sequence $\hat{\beta}$ which is strictly decreasing such that $E(\hat{\beta})= \pm E(\beta)$, with the sign depending on the number of order reversals. We observe that if $\hat{\beta}_{1}>n$, then $E(\hat{\beta})=0$. Furthermore, if $\hat{\beta}_{1}=n$, then the initial $e_{n}^{n}$ is the only nonzero entry in the first row of the determinant defining $E(\hat{\beta})$ and so we may express it as $e_{n}^{n}$ times the lower right minor. In this way, we see that

$$
E(\hat{\beta})=\left(e_{n}^{n}\right)^{k} E\left(\beta^{\circ}\right)
$$

for some $k \in \mathbb{Z}_{\geq 0}$ and $\beta^{\circ}$ a strictly decreasing sequence of integers, such that $\beta_{m}^{\circ}=1-m$ for $m$ sufficiently large and $\beta_{1}^{\circ}<n$. We may then recover two partitions $\lambda^{\circ}$ and $\mu^{\circ}$ such that

$$
\left(\lambda^{\circ}\right)_{1}^{\prime}+\left(\mu^{\circ}\right)_{1}^{\prime} \leq n, \quad \mu_{1}^{\circ}=s-k
$$

and

$$
\beta\left(a\left(\lambda^{\circ}, \mu^{\circ}\right)\right)=\beta^{\circ}
$$

Thus we have

$$
\chi_{G L_{n}}(\bar{\mu} ; \lambda)= \pm \chi_{G L_{n}}\left(\bar{\mu}^{\circ} ; \lambda^{\circ}\right)
$$

with the sign being the one obtained from the order reversals to get $\hat{\beta}$ from $\beta$.
11.6. Dictionary. Now, for the reader familiar with [B-K-W] or interested in consulting it now, we offer a partial dictionary to translate from our notation to theirs.

The translation is not precise since the notion of a projective limit does not appear in [B-K-W], only various identities which "stabilize" once $n$ is sufficiently large and hence are ripe for interpretation in terms of this notion. In the text of this appendix, we have used $\lambda$ and $\mu$ both for partitions (which might not be weights) and for weights (which might not be partitions); throughout this table, $\lambda$ and $\mu$ are partitions, while other Greek letters are weights. Some of the notations depend on the number of parts of the partitions. These will be denoted by $p$ and $q$ respectively. We subdivide into three parts: notations common in the majority of $[\mathrm{B}-\mathrm{K}-\mathrm{W}]$, notations appearing mainly in section 6 , and then a few logical equivalencies. We remark that $[\lambda]_{ \pm}$and $[\Delta ; \lambda]_{ \pm}$are also defined for $p>n$. For a "universal character" interpretation, see $[\mathrm{Ko}]$. Note that $[\bar{\emptyset} ; \lambda]_{ \pm}$and $[\Delta ; \bar{\emptyset} ; \lambda]_{ \pm}$are not equal to $[\lambda]_{ \pm}$and $[\Delta ; \lambda]_{ \pm}$for $p>n$.

| Our notation | Their notation |
| :---: | :---: |
| det | $\varepsilon$ |
| $L R_{\lambda, \mu}^{\nu}$ | $m_{\lambda \mu}^{\nu}$ |
| $\pi_{n} \circ \rho_{n}\left(\chi_{G L}(\bar{\mu} ; \lambda)\right)$ | $\{\bar{\mu} ; \lambda\}$ |
| $\tilde{\rho}_{n}\left(\chi_{G L}(\lambda)\right)=\pi_{n} \circ \rho_{n}\left(\chi_{G L}(\bar{\emptyset} ; \lambda)\right)$ | $\{\lambda\}$ |
| $\pi_{n} \circ \rho_{n}\left(\chi_{G L}(\bar{\mu} ; \emptyset)\right)$ | $\{\bar{\mu}\}$ |
| $\pi_{S p_{2 n}}\left(\chi_{S p}(\lambda)\right)$ | $\langle\lambda\rangle$ |
| $\iota \quad{ }^{\text {chen }}$ | $\dagger$ or $\left.\cdots]_{+} \rightarrow \cdots\right]_{-}$ |
| $\chi_{S O_{2 n}}(\lambda)$ | $[\lambda]_{+} \quad$ if $p \leq n$ |
| $\chi_{\text {SO }}^{2 n}$ ( $\left.\lambda\right)+\chi_{\text {SO }}^{2 n}$ ( $\lambda^{\iota}$ ) | $[\lambda] \quad$ if $p=n$ |
|  | $[\lambda]=[\lambda]_{+}=[\lambda]_{-} \quad$ if $p<n$ |
| $\chi_{\text {Spin }_{2 n}}\left(\lambda+\left(\frac{1}{2}\right)^{n}\right)$ | $[\Delta ; \lambda]_{+} \quad$ if $p \leq n$ |
| $\chi_{\text {Spin }_{2 n}}\left(\lambda+\left(\frac{1}{2}\right)^{n}\right)+\chi_{\text {SO }}^{2 n}$ ( $\left.\left(\lambda+\left(\frac{1}{2}\right)^{n}\right)^{\iota}\right)$ | $[\Delta ; \lambda] \quad$ if $p \leq n$ |
| $\chi_{\operatorname{Spin}_{2 n}}\left(\lambda+1^{n}\right)$ | $[\square ; \lambda]_{+} \quad$ if $p \leq n$ |
| $\chi_{{\text {Spin } 2_{2 n}}\left(\lambda+1^{n}\right)+\chi_{\text {Spin }_{2 n}}\left(\left(\lambda+1^{n}\right)^{\iota}\right)}$ | $[\square ; \lambda] \quad$ if $p \leq n$ |
| $T_{\widetilde{G L}_{n} \rightarrow \operatorname{Spin}_{2 n}} \circ \pi_{n} \circ \rho_{n}\left(\chi_{G L}(\bar{\mu} ; \lambda)\right)$ | $[\bar{\mu} ; \lambda]_{+}$ |
| $T_{\widetilde{G L}_{n} \rightarrow \text { Spin }_{2 n}}\left(\operatorname{det}^{\frac{1}{2}} \pi_{n} \circ \rho_{n}\left(\chi_{G L}(\bar{\mu} ; \lambda)\right)\right)$ | $[\Delta ; \bar{\mu} ; \lambda]_{+}$ |
| $\chi_{G}(\nu)$ | $(\nu)_{G}$ (alternate notation of section 6.) |
| $\rho_{G}$ |  |
| $c(H)_{\sigma, \nu+\rho_{G}-\rho_{H}}^{\tau}$ | $K_{\sigma_{H}, \nu_{G}+\delta_{G}-\delta_{H}}^{\tau_{H}}$ |
| $n_{\eta}^{\nu}$ | $B_{\nu_{G}}^{\eta_{H}, \nu_{G}}{ }^{\text {a }}$ |
| $m_{\nu}^{G}(\sigma)$ | $M_{\nu_{G}}^{\sigma_{T}}$ |
| Equation (11.4.3) | Equations (7.1a) \& (7.1b) |
| Equation (11.4.4) | Equations (7.5a) \& (7.5b) |
| $(\bar{\mu} ; \lambda) \mapsto \pm\left(\bar{\mu}^{\circ} ; \lambda^{\circ}\right)$ | $U_{n}, S U_{n}$ Modification rule of Table 3 (possibly applied more than once) |
| $\chi_{G}(\nu) \mapsto \frac{\nu}{\|n u\|_{G}} \chi_{G}\left(\|\nu\|_{G}\right)$ | Modification rule for $G$ in Tables $4 \& 5$ (possibly applied more than once) |

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