ON THE DISTRIBUTION OF SATAKE PARAMETERS for Siegel Modular Forms

Andrew Knightly and Charles Li

Received: October 19, 2017 Revised: February 5, 2019

Communicated by Don Blasius

ABSTRACT. We prove a harmonically weighted equidistribution result for the p-th Satake parameters of the family of automorphic cuspidal representations of $PGSp(2n)$ of fixed weight k and prime-to-p level $N \to \infty$. The main tool is a new asymptotic Petersson formula for $PGSp(2n)$ in the level aspect.

2010 Mathematics Subject Classification: 11F46 (11F25, 11F30, 11F70, 11F72)

Keywords and Phrases: Siegel modular forms, Petersson formula, Satake parameters

CONTENTS

678 Andrew Knightly and Charles Li

1 Introduction

For a split reductive algebraic group G over a number field F , let $\mathcal{A}(G)$ denote the set of cuspidal automorphic representations of $G(\mathbf{A}_F)$. Each element of $\mathcal{A}(G)$ factorizes as a restricted tensor product $\pi = \otimes_v \pi_v$ of irreducible representations of the local groups $G(F_v)$. If v is a nonarchimedean place of F, then the unramified irreducible representations of $G(F_v)$ are parametrized (via the Satake isomorphism) by the semisimple conjugacy classes in the complex dual group $\widehat{G} = \widehat{G}(\mathbf{C})$. When π_v is unramified, we let

$$
t_{\pi_v} \in \widehat{T}/W
$$

denote the associated Satake parameter. Here, $\hat{T} = \hat{T}(\mathbf{C})$ is a split maximal torus of \widehat{G} and $W = N_G(T)/T$ is the Weyl group of G.

It is of great interest to understand the distribution of the points t_{π_v} , possibly with weights, as π and/or v vary. If π is fixed and v varies, then according to the general Sato-Tate conjecture, the points are expected to be equidistributed relative to some naturally defined measure μ_{ST} with support in the maximal compact subgroup of \hat{T} . The support condition reflects the presumed predominance of representations satisfying the Ramanujan conjecture.

Many people have considered the analogous "vertical" question of fixing v and varying π in a family, starting with Birch's 1968 investigation of the problem for elliptic curves over $\mathbf{Z}/p\mathbf{Z}$, [\[Bi\]](#page-67-0). A table summarizing the numerous subsequent vertical equidistribution results for GL(2) families is given in [\[KR\]](#page-68-0). Extending work of Serre and others on the case of classical holomorphic cusp forms, Shin proved generally that if G admits discrete series over **R**, then when the t_{π_v} are unweighted, the relevant measure (for many natural families) is the Plancherel measure at v, [\[Shin\]](#page-69-0). Shin and Templier also obtained a quantitative version of this result with error bounds, [\[ST\]](#page-70-0). Applications include: (i) a diagonal hybrid where the size of the family and the place v both tend to infinity (the relevant measure being Sato-Tate rather than Plancherel), and (ii) a determination of the distribution of the low-lying zeroes of certain families of automorphic Lfunctions for G. Matz and Templier have recently treated the case of $GL(n)$. [\[MT\]](#page-69-1). We refer to Sarnak, Shin and Templier [\[SST\]](#page-69-2) for a precise formulation of various Sato-Tate problems and related topics.

When the Satake parameters at a fixed place v are given the harmonic weights that arise naturally in the Petersson/Kuznetsov trace formula, it has been shown for many $GL(2)$ and $GL(3)$ families that they exhibit equidistribution relative to the Sato-Tate measure itself rather than the Plancherel measure ([\[Br\]](#page-67-1), [\[Li\]](#page-69-3), [\[KL3\]](#page-68-1), [\[BrM\]](#page-67-2), [\[BBR\]](#page-67-3), [\[Z\]](#page-70-1)). In this paper we consider the distribution of harmonically weighted Satake parameters at a fixed place p for the group $G = GSp(2n)$. For simplicity we work over **Q** and assume trivial central character. We consider cuspidal representations π of level N for which π_{∞} is a fixed holomorphic discrete series representation of weight $k > 2n$. We weight each Satake parameter t_{π_p} by the globally defined value

$$
w_{\pi} = \sum_{\varphi \in E_{\mathbf{k}}(\pi)} \frac{|c_{\sigma}(\varphi)|^2}{\|\varphi\|^2},\tag{1.1}
$$

where $E_{\mathbf{k}}(\pi)$ is a finite orthogonal set of cuspidal Hecke eigenforms giving rise to π , and $c_{\sigma}(\varphi)$ denotes a Fourier coefficient, defined in [\(6.2\)](#page-21-1). Weighted in this way, we prove that the parameters become equidistributed relative to a certain probability measure μ as $N \to \infty$ (see Theorem [8.3\)](#page-34-0). In contrast to the $GL(2)$ case, the measure depends on p . Subject to a natural hypothesis on the growth of the geometric side of the trace formula (which holds at least when $n = 2$, in Theorem [8.4](#page-36-1) we relate μ to the Sato-Tate measure. In particular, the

hypothesis implies that μ is supported on the tempered spectrum, and tends to the Sato-Tate measure as $p \to \infty$.

A very similar equidistribution problem has been studied already in the case $n = 2$ $n = 2$ $n = 2$ by Kowalski, Saha, and Tsimerman ([\[KST2\]](#page-68-2)), who fix the level $N = 1$ and let the archimedean parameter $k \to \infty$. Using a formula of Sugano, they were able to form a connection between Satake parameters and Fourier coefficients, the latter of which they control with the intricate Petersson formula for GSp(4) due to Kitaoka [\[Ki\]](#page-68-3). Their weights involve certain linear combinations of Fourier coefficients and depend on the choice of a class group character. In a special case they coincide with the ones given in [\(1.1\)](#page-2-0) with $\sigma = I_n$. Their methods have been adapted to treat the case of higher level $N \to \infty$ by M. Dickson [\[D\]](#page-67-4). Kim, Wakatsuki and Yamauchi have also investigated the equidistribution problem for $GSp(4)$, via Arthur's invariant trace formula [\[KWY\]](#page-68-4). In all of these works, a quantitative equidistribution statement is proven, with application to the distribution of low-lying zeros of L-functions.

The key technical tools used in [\[KST2\]](#page-68-2), namely, Kitaoka's formula and Sugano's formula, are not yet available when $n > 2$. Nevertheless, we can apply two simple ideas to treat the higher rank case. The first is to use a Hecke operator as a test function in the relative trace formula to derive a Petersson formula for $GSp(2n)$ whose spectral side involves two Fourier coefficients (as usual) with the additional inclusion of a Satake parameter. In this way we can access the Satake parameters directly without the use of Sugano's formula. In order to project onto the holomorphic cusp forms of weight k, we use a certain matrix coefficient of the weight k holomorphic discrete series π_k of $GSp(2n, \mathbf{R})$ as the archimedean component of our test function. This function is computed explicitly in the Appendix (Theorem [A.9\)](#page-61-0).

The second idea is to take the limit of the kernel function before integrating. This allows us to avoid computing or estimating all but a few of the orbital integrals that show up on the geometric side. The result is an asymptotic Petersson formula with only a finite sum on the geometric side (Theorem 6.3).^{[†](#page-3-1)} In the simplest case where the Hecke operator is trivial, it is given as follows (see §[7.3,](#page-27-1) where notation is explained in detail).

THEOREM 1.1. Let $\mathcal{B}_k(N)$ be an orthogonal basis of the space of degree n Siegel cusp forms of weight $k > 2n$ with $2|nk$, and level group $\Gamma_0(N) = \{ \begin{pmatrix} A & B \\ C & D \end{pmatrix} \in$ $\operatorname{Sp}_{2n}(\mathbf{Z})|C \equiv O \mod N$. For symmetric positive-definite half-integral matrices σ_1, σ_2 ,

$$
\lim_{N \to \infty} \frac{1}{\psi(N)} \sum_{F \in \mathcal{B}_k(N)} \frac{a_{\sigma_1}(F) \overline{a_{\sigma_2}(F)}}{\|F\|^2} = \delta_k(\sigma_1, \sigma_2) c_{n \mathbf{k} \sigma_1},\tag{1.2}
$$

[†]After this paper was written, we became aware of [\[KST1\]](#page-68-5), in which a similar idea is applied to classical Siegel Poincaré series to obtain an analog of Theorem [1.1](#page-3-0) for $k \to \infty$.

where $a_{\sigma_j}(F)$ are Fourier coefficients, $\psi(N) = [\Gamma_0(1) : \Gamma_0(N)],$

$$
c_{n\mathbf{k}\sigma_1} = \frac{(\det \sigma_1)^{\mathbf{k} - (n+1)/2}}{\pi^{n(n-1)/4} (4\pi)^{n(n+1)/2 - n\mathbf{k}} \prod_{j=1}^n \Gamma(\mathbf{k} - \frac{n+j}{2})},
$$

and

$$
\delta_{\mathbf{k}}(\sigma_1, \sigma_2) = \sum_{\substack{A \in \mathrm{GL}_n(\mathbf{Z}) / \{\pm I_n\} \\ t_{A\sigma_1 A = \sigma_2}}} \det(A)^{\mathbf{k}}.
$$
\n(1.3)

Remarks:

(1) When k is even, $\delta_{\mathbf{k}}(\sigma_1, \sigma_2) = \delta_0(\sigma_1, \sigma_2) = \#\{A \in \mathrm{GL}_n(\mathbf{Z})/\{\pm I_n\}\}\,{}^t\!A\sigma_1A =$ σ_2 . When k is odd, it may happen that $\delta_k(\sigma_1, \sigma_2) = 0$ even when $\delta_0(\sigma_1, \sigma_2) > 0$ (for example, if σ_1, σ_2 are diagonal). In such cases, the spectral side also vanishes, since $a_{\sigma_j}(F) = 0$ by [\[Kl,](#page-68-6) p. 45]. We note that $\delta_{\mathbf{k}}(\sigma_1, \sigma_2)$ does not always vanish when k is odd.

(2) When $n = 2$ and k is even, the above result is shown in [\[CKM\]](#page-67-5), [\[D\]](#page-67-4). (In [\[CKM,](#page-67-5) Theorem 1.1, Remark 1.4], it is incorrectly asserted that [\(1.2\)](#page-3-2) holds when $n = 2$ and k is odd, but with $\delta_0(\sigma_1, \sigma_2)$ in place of $\delta_k(\sigma_1, \sigma_2)$. This is incompatible with the observations in our first remark.)

(3) $c_{n\mathbf{k}\sigma_1}$ is the constant of proportionality between $a_{\sigma_1}(F)$ and a suitably normalized inner product of F with the σ_1 -Poincaré series of weight k ([\[M\]](#page-69-4), [\[Kl,](#page-68-6) p. 90], [\[D,](#page-67-4) Lemma 6.2]).

(4) Although we have highlighted the above special case, the main focus of this paper is on local Satake parameters, which are absent in the above theorem.

Petersson/Kuznetsov trace formulas play a fundamental role in the study of automorphic L-functions. There are well-established methods for $GL(2)$, and to a lesser extent $GL(3)$, but applications to other groups are rare. For a recent example, Blomer has used Kitaoka's formula to compute first and second moments of spinor L -functions of $GSp(4)$ of full level and large weight, with power-saving error term, [\[Bl\]](#page-67-6). Waibel subsequently treated the case of large prime level, [\[W\]](#page-70-2). This requires a close analysis of the generalized Kloosterman sums appearing in the off-diagonal terms. It would be of great interest to extend Kitaoka's formula from degree 2 to degree n . The machinery we develop here can form the starting point for such a generalization. Although at present there is no quantification of the error term for finite N if $n > 2$ (the $n = 2$ case is treated in Appendix [B\)](#page-63-0), the asymptotic formula is sufficient for obtaining the equidistribution result.

A natural question is whether the same method can be used to obtain the weighted equidistribution of Satake parameters in the case of $k \to \infty$. Although we certainly expect such a result to hold, it does not seem possible to show it by the above method without a more detailed consideration of the off-diagonal terms. This is discussed in a remark after the proof of Proposition [6.1.](#page-21-2)

We close this introduction with some additional comments comparing our re-sults and methods to those in [\[KST2\]](#page-68-2). If F is a Siegel cusp form of degree $n = 2$ which is unramified at p, it determines an unramified local representation π_p of GSp(4). Sugano's formula expresses the value of a Laurent polynomial $U_p^{l,m}$ evaluated at the Satake parameters of π_p as a value of the spherical function in the Bessel model for π_p . The latter can in turn be expressed as a certain linear combination of Fourier coefficients of F , with coefficients involving a character of a class group. This is the source of the weights in the equidistri-bution result of [\[KST2\]](#page-68-2). This method has the advantage that the $U_p^{l,m}$ were shown by Furusawa and Shalika to be orthonormal with respect to a certain explicit local measure η_p on the space \mathfrak{X}_p (defined near [\(8.1\)](#page-33-2) below) containing the Satake parameters. The latter fact leads directly to the weighted equidistribution result, since a refinement of Kitaoka's formula has the main term $\delta(l,m) = \langle U_p^{l,m}, U_p^{0,0} \rangle = \int_{\mathfrak{X}_p} U_p^{l,m} d\eta_p$. By contrast, the measure μ in our Theorem [8.3](#page-34-0) is not given explicitly, because we have not found an orthonormal class of functions for which the main term of the trace formula has such a simple form.

Acknowledgements: We would like to thank the referee, whose detailed comments led to significant improvements in the exposition. We thank Yuk-Kam Lau, Nigel Pitt, Abhishek Saha, and Nicolas Templier for helpful discussions. We are also grateful for financial support from the University of Maine Office of the Vice President for Research and from the Simons Foundation. The first author would also like to thank the Department of Mathematics at the Chinese University of Hong Kong for its hospitality.

This work was partially supported by a grant from the Simons Foundation $(\#317659 \text{ to the first author}).$

2 The Satake transform

Here we recall some basic background about the Satake transform. References include $\operatorname{[Gr]}$ and $\operatorname{[Sha]}$. For notation in this section, let G be a split group defined over \mathbf{Q}_p , and let T be a (split) maximal torus of G defined over \mathbf{Q}_p and contained in a Borel subgroup $B = TN$ with N unipotent. Let $X^*(T) =$ Hom(T, GL₁) denote the lattice of algebraic characters of T, and $X_*(T) =$ $Hom(GL_1, T)$ the cocharacter lattice.

For each prime p, let $K_p = G(\mathbf{Z}_p)$, which is a maximal compact subgroup of $G_p = G(\mathbf{Q}_p)$. Let $\mathcal{H}(G_p, K_p)$ be the Hecke algebra of locally constant compactly supported complex-valued bi- K_p -invariant functions on G_p . The Satake transform of an element $f \in \mathcal{H}(G_p, K_p)$ is the function on $T(\mathbf{Z}_p)\setminus T(\mathbf{Q}_p)$ given by

$$
\mathcal{S}f(t) = \delta(t)^{1/2} \int_{N(\mathbf{Q}_p)} f(tn)dn.
$$
 (2.1)

Here $\delta(t) = |\det \text{Ad}(t)|_{\text{Lie}(N_p)}|$ is the modular function for B_p ^{\ddagger} By way of motivation for (2.1) , suppose

$$
\pi_{\chi} = \mathrm{Ind}_{B_p}^{G_p}(\chi\delta^{1/2})
$$

is an unramified representation of G_p , with nonzero K_p -invariant vector ϕ . Then for $f \in \mathcal{H}(G_p, K_p)$, ϕ is an eigenfunction of the operator $\pi_\chi(f)$, with eigenvalue

$$
\omega_{\chi}(f) = \int_{T(\mathbf{Q}_p)} \mathcal{S}f(t)\chi(t)dt.
$$
\n(2.2)

The Satake transform is a C-algebra isomorphism

$$
S: \mathcal{H}(G_p, K_p) \longrightarrow \mathcal{H}(T(\mathbf{Q}_p), T(\mathbf{Z}_p))^W, \tag{2.3}
$$

where the latter denotes the elements which are fixed by the Weyl group W. Let $\hat{T} = \text{Hom}(X_*(T), \mathbf{C})$ be the dual group of T. It satisfies

$$
X^*(T) \cong X_*(\widehat{T}), \quad X_*(T) \cong X^*(\widehat{T}).
$$

Using the algebra isomorphisms

$$
\mathcal{H}(T(\mathbf{Q}_p), T(\mathbf{Z}_p))^W \cong \mathbf{C}[X_*(T)]^W \cong \mathbf{C}[X^*(\widehat{T})]^W,
$$
\n(2.4)

we may view Sf as a function on any of these three spaces. We explain this in some more detail. The first isomorphism in [\(2.4\)](#page-6-1) arises by identifying an element of $T(Q_p)/T(Z_p)$ with a tuple of integer powers of p, which is of the form $\lambda(p)$ for a unique $\lambda \in X_*(T)$. Thus if we write

$$
Sf = \sum_{t \in T(\mathbf{Q}_p)/T(\mathbf{Z}_p)} a_t C_t \in \mathcal{H}(T(\mathbf{Q}_p), T(\mathbf{Z}_p))
$$
(2.5)

where $a_t \in \mathbf{C}$ is nonzero for at most finitely many t, and C_t is the characteristic function of the coset t , we can make the identification

$$
\mathcal{S}f = \sum_{\lambda \in X_*(T)} a_{\lambda} \lambda \in \mathbf{C}[X_*(T)],
$$

where $a_{\lambda} = a_{\lambda(p)}$ from [\(2.5\)](#page-6-2). Fix an isomorphism $X_*(T) \cong X^*(T)$ and denote it by $\lambda \mapsto \hat{\lambda}$. Then we may in turn identify Sf with the function

$$
Sf = \sum_{\hat{\lambda} \in X^*(\hat{T})} a_{\hat{\lambda}} \hat{\lambda} \in \mathbf{C}[X^*(\hat{T})],\tag{2.6}
$$

where $a_{\hat{\lambda}} = a_{\lambda}$.

[‡]Later on we will take $G = \text{GSp}(2n)$ and B the Borel subgroup determined by the set of positive roots chosen in §[3.2.](#page-9-0) Then $\delta(t)^{1/2} = p^{-\langle \lambda, \rho \rangle}$ if $t = \lambda(p)$ for $\lambda \in X_*(T)$ and ρ is given by (3.10) .

An unramified character χ of $T(Q_p)$ as in [\(2.2\)](#page-6-3) can be identified with the Satake parameter $t_{\chi} \in \widehat{T}(\mathbf{C})$ determined by

$$
\chi(\lambda(p)) = \hat{\lambda}(t_{\chi}) \quad \text{for all } \lambda \in X_*(T)
$$
\n(2.7)

(cf. $[Ca, p. 134, Eq. (3)]$). With this notation, using (2.5) and (2.6) , and taking meas $(T(\mathbf{Z}_p)) = 1$, [\(2.2\)](#page-6-3) becomes

$$
\omega_{\chi}(f) = \sum_{\lambda} a_{\lambda} \chi(\lambda(p)) = \sum_{\hat{\lambda}} a_{\hat{\lambda}} \hat{\lambda}(t_{\chi}) = \mathcal{S}f(t_{\chi}). \tag{2.8}
$$

When $\pi = \pi_\chi$ is given, we write $t_\pi = t_\chi$ for the Satake parameter of π .

PROPOSITION 2.1. Viewing Sf as a function on \hat{T} as in [\(2.6\)](#page-6-4), we have

$$
\overline{\mathcal{S}f}=\mathcal{S}f^*,
$$

where $f^*(g) = \overline{f(g^{-1})}.$

Proof. For $t \in T$, we need to show that $Sf(t) = Sf^*(t)$. As in [\(2.7\)](#page-7-2), there exists a unique unramified character χ of $T(Q_p)$ such that $\chi(\lambda(p)) = \hat{\lambda}(t)$ for all $\lambda \in X_*(T)$. Recalling that $\pi_\chi(f)^* = \pi_\chi(f^*)$, by [\(2.8\)](#page-7-3) we have, for the spherical unit vector $\phi \in \pi_{\chi}$,

$$
\mathcal{S}f(t) = \omega_{\chi}(f) = \langle \pi_{\chi}(f)\phi, \phi \rangle = \langle \phi, \pi_{\chi}(f^*)\phi \rangle = \overline{\omega_{\chi}(f^*)} = \overline{\mathcal{S}f^*(t)}.
$$

3 The symplectic group

3.1 DEFINITION

Henceforth, we will denote by G the algebraic group GSp_{2n} defined as follows. For any commutative ring R, let $M_{2n}(R)$ be the set of $2n \times 2n$ matrices with entries in R. Letting O denote the zero-matrix of suitable dimension, and I_n denote the $n \times n$ identity matrix, define

$$
J = \begin{pmatrix} O & I_n \\ -I_n & O \end{pmatrix},
$$

$$
Sp_{2n}(R) = \{ M \in M_{2n}(R) | {}^{t}MJM = J \},
$$

 $\text{GSp}_{2n}(R) = \{ M \in M_{2n}(R) | {}^{t}MJM = r(M)J \text{ and } r(M) \text{ is a unit in } R \}.$

Thus a matrix $M = \begin{pmatrix} A & B \\ C & D \end{pmatrix}$ belongs to $GSp_{2n}(R)$ if and only if there exists a similitude $r(M) \in R^*$ such that

$$
{}^{t}AC = {}^{t}CA, \quad {}^{t}BD = {}^{t}DB, \quad {}^{t}AD - {}^{t}CB = r(M)I_n. \tag{3.1}
$$

Taking inverses in ${}^tMJM = r(M)J$ shows that ${}^tM = \begin{pmatrix} {}^t A & {}^t C \\ {}^t B & {}^t D \end{pmatrix} \in \text{GSp}_{2n}(R)$, and the above relations applied to this matrix give

$$
A^t B = B^t A
$$
, $C^t D = D^t C$, $A^t D - B^t C = r(M) I_n$. (3.2)

Note that $\text{Sp}_{2n}(R) = \{M \in \text{GSp}_{2n}(R) | r(M) = 1\}.$ Define

$$
\mathbb{G}=\mathrm{PGSp}_{2n}=\mathrm{GSp}_{2n}\, /Z,
$$

where Z is the center (the set of scalar matrices).

Let $S_n(R)$ be the set of $n \times n$ symmetric matrices over R. The Siegel upper half space is the following set of complex symmetric matrices

$$
\mathcal{H}_n = \{ X + iY \in S_n(\mathbf{C}) | X, Y \in S_n(\mathbf{R}), Y > 0 \},\tag{3.3}
$$

where $Y > 0$ means that Y is positive definite. It is a complex vector space of dimension $\frac{n(n+1)}{2}$. Letting

$$
\operatorname{GSp}_{2n}(\mathbf{R})^+ = \{ M \in \operatorname{GSp}_{2n}(\mathbf{R}) | r(M) > 0 \},\
$$

there is a transitive action of $\text{GSp}_{2n}(\mathbf{R})^+$ on \mathcal{H}_n given by

$$
M \cdot 3 = (A3 + B)(C3 + D)^{-1}
$$
 $(3 \in \mathcal{H}_n).$

The stabilizer in $\text{Sp}_{2n}(\mathbf{R})$ of the element $iI_n \in \mathcal{H}_n$ is the compact subgroup

$$
K_{\infty} = \left\{ \begin{pmatrix} A & B \\ -B & A \end{pmatrix} \in \text{Sp}_{2n}(\mathbf{R}) | A + iB \in \text{U}(n) \right\},\tag{3.4}
$$

where U(n) is the group of $n \times n$ complex unitary matrices X (so $X^{-1} = {}^t \overline{X}$). Define the Siegel parabolic subgroup

$$
P(\mathbf{R}) = \left\{ \begin{pmatrix} A & B \\ O & D \end{pmatrix} \in \text{GSp}_{2n}(\mathbf{R}) \right\}
$$

=
$$
\left\{ \begin{pmatrix} A & B \\ r & 1 \end{pmatrix} \begin{pmatrix} I & S \\ I & I \end{pmatrix} | A \in \text{GL}_n(\mathbf{R}), r \in \mathbf{R}^*, S \in S_n(\mathbf{R}) \right\}
$$
(3.5)

(where O denotes the $n \times n$ zero matrix), and set $\mathbb{P} = P/Z$. We recall the decomposition

$$
\mathrm{GSp}_{2n}(\mathbf{R}) = P(\mathbf{R})K_{\infty}.
$$

For a prime p , let

$$
K_p = \mathrm{PGSp}_{2n}(\mathbf{Z}_p).
$$

Haar measure on $PGSp_{2n}(\mathbf{Q}_p)$ will be normalized so that $meas(K_p) = 1$. For an integer $N > 0$, set

$$
K_0(N)_p = \left\{ \begin{pmatrix} A & B \\ C & D \end{pmatrix} \in K_p | C \equiv O \mod N \mathbf{Z}_p \right\},\tag{3.6}
$$
\n
$$
K_0(N) = \prod_{p < \infty} K_0(N)_p.
$$

3.2 ROOT DATA

We review some standard material to fix notation and terminology that will be used in the sequel. Let F be an algebraically closed field. In this subsection we write G for $G(F)$, and similarly for the other algebraic groups considered. In G, the diagonal subgroup

$$
T = \{ t = \text{diag}(a_1, \dots, a_n, \frac{r}{a_1}, \dots, \frac{r}{a_n}) | ra_1 \cdots a_n \neq 0 \}
$$
 (3.7)

is a split maximal torus. Given $\chi \in X^*(T)$, there exist $k_0, \ldots, k_n \in \mathbb{Z}$ such that

$$
\chi(t) = r^{k_0} a_1^{k_1} \cdots a_n^{k_n}
$$
\n(3.8)

for t as in [\(3.7\)](#page-9-2). By associating χ with the tuple (k_0, \ldots, k_n) , we obtain a natural identification $X^*(T) \cong \mathbb{Z}^{n+1}$. Let $e_j \in \mathbb{Z}^{n+1}$ be the $(j+1)$ -th standard basis vector $(j = 0, \ldots, n)$. The set Φ of roots (for the action of T on Lie(G)) consists of \pm (e) \pm 2 $(1 \times i \times i \times n)$

$$
\pm (e_j - e_i) \quad (1 \le i < j \le n),
$$
\n
$$
\pm (e_0 - e_j - e_i) \quad (1 \le i \le j \le n). \tag{3.9}
$$

We take the set Φ^+ of positive roots to consist of those which have the $+$ coefficient. The corresponding set of simple roots is

$$
\Delta = \{e_{j+1} - e_j | j = 1, \ldots, n-1\} \cup \{e_0 - 2e_n\}.
$$

Let

$$
\rho = \frac{1}{2} \sum_{\chi \in \Phi^+} \chi = \frac{1}{2} \left(\sum_{1 \le i < j \le n} (e_j - e_i) + \sum_{1 \le i \le j \le n} (e_0 - e_i - e_j) \right).
$$

Explicitly,

$$
\rho = \frac{n(n+1)}{4}e_0 - ne_1 - (n-1)e_2 - \dots - e_n.
$$
\n(3.10)

The cocharacter lattice is $X_*(T) = \text{Hom}(F^*,T)$. We identify a tuple $\lambda =$ $(\ell_0, \ell_1, \ldots, \ell_n) \in \mathbf{Z}^{n+1}$ with the cocharacter

$$
\lambda(a) = \text{diag}(a^{\ell_1}, \dots, a^{\ell_n}, a^{\ell_0 - \ell_1}, \dots, a^{\ell_0 - \ell_n}).
$$
\n(3.11)

In this way, $X_*(T) \cong \mathbf{Z}^{n+1}$.

The composition of a character with a cocharacter yields a rational homomorphism $F^* \to F^*$, which is necessarily of the form $x \mapsto x^m$ for $m \in \mathbb{Z}$. Thus, we have a natural pairing $X^*(T) \times X_*(T) \to \mathbb{Z}$ given by

$$
\chi(\lambda(x)) = x^{\langle \chi, \lambda \rangle}.\tag{3.12}
$$

In terms of the coordinates given above, this works out to

$$
\langle (k_0,\ldots,k_n), (\ell_0,\ldots,\ell_n) \rangle = k_0 \ell_0 + \cdots + k_n \ell_n. \tag{3.13}
$$

Our main group of interest is $\mathbb{G} = \text{PGSp}_{2n}$. Letting $\mathbb{T} = T/Z$ be the maximal torus, we can identify its character lattice $X^*(\mathbb{T})$ with the subset of $X^*(T)$ consisting of all characters which are trivial on Z. Thus,

$$
X^*(\mathbb{T}) \cong \{(k_0, k_1, \dots, k_n) \in \mathbb{Z}^{n+1} | 2k_0 + k_1 + \dots + k_n = 0\}.
$$
 (3.14)

Note that the roots [\(3.9\)](#page-9-3) belong to this set, so we may identify Φ with the set of roots in $X^*(\mathbb{T})$. The cocharacter lattice $X_*(\mathbb{T})$ can be viewed as the quotient of $X_*(T)$ by the subgroup of cocharacters taking values in Z. In terms of the coordinates above [\(3.11\)](#page-9-4), we have

$$
X_*(\mathbb{T}) \cong \mathbf{Z}^{n+1}/(2, 1, \dots, 1)\mathbf{Z}.\tag{3.15}
$$

The pairing [\(3.12\)](#page-9-5) makes sense for $(\chi, \lambda) \in X^*(\mathbb{T}) \times X_*(\mathbb{T})$, and the formula (3.13) is independent of the choice of coset representative for λ .

The Weyl group of G, namely

$$
W = N_{\mathbb{G}}(\mathbb{T})/Z_{\mathbb{G}}(\mathbb{T}) = N_G(T)/Z_G(T),
$$

acts on T (and also T) by conjugation. It is isomorphic to $S_n \ltimes (\mathbf{Z}/2\mathbf{Z})^n$, with the following generators:

$$
t \mapsto \text{diag}(a_{\sigma(1)}, \ldots, a_{\sigma(n)}, \frac{r}{a_{\sigma(1)}}, \ldots, \frac{r}{a_{\sigma(n)}})
$$

for σ in the symmetric group S_n , and, for $1 \leq i \leq n$,

$$
t \mapsto \text{diag}(a_1,\ldots,a_{i-1},\frac{r}{a_i},a_{i+1},\ldots,a_n,\frac{r}{a_1},\ldots,\frac{r}{a_{i-1}},a_i,\frac{r}{a_{i+1}},\ldots,\frac{r}{a_n}).
$$

Likewise, W acts faithfully on $X^*(\mathbb{T})$ by

$$
w\chi(t) = \chi(wtw^{-1}).
$$

The corresponding generators are

$$
(k_0, k_1, \dots, k_n) \mapsto (k_0, k_{\sigma^{-1}(1)}, \dots, k_{\sigma^{-1}(n)})
$$
\n(3.16)

and

$$
(k_0, k_1, \ldots, k_n) \mapsto (k_0 + k_i, k_1, \ldots, k_{i-1}, -k_i, k_{i+1}, \ldots, k_n). \tag{3.17}
$$

Using the pairing [\(3.12\)](#page-9-5), an action of W on $X_*(\mathbb{T})$ is defined implicitly via

$$
\langle w\chi, w\lambda\rangle = \langle \chi, \lambda\rangle.
$$

In terms of the ℓ -coordinates in [\(3.11\)](#page-9-4), the action on $X_*(T)$ of the Weyl element in (3.16) is given by

$$
(\ell_0, \ell_1, \ldots, \ell_n) \mapsto (\ell_0, \ell_{\sigma^{-1}(1)}, \ldots, \ell_{\sigma^{-1}(n)}),
$$

and the one corresponding to [\(3.17\)](#page-10-1) is given by

$$
(\ell_0,\ell_1,\ldots,\ell_n)\mapsto (\ell_0,\ell_1,\ldots,\ell_{i-1},\ell_0-\ell_i,\ell_{i+1},\ldots,\ell_n).
$$

Suppose $\chi, \chi' \in X^*(\mathbb{T})$ correspond respectively to (k_i) , $(k'_i) \in \mathbb{Z}^{n+1}$ as in [\(3.14\)](#page-10-2). Then the pairing

$$
(\chi,\chi')=\sum_{k=1}^n k_i k_i'
$$

is W -invariant. (This is easily verified using the relation in (3.14) .) For a root $\alpha \in \Phi$, there is a unique coroot $\alpha^{\vee} \in X_*(\mathbb{T})$ satisfying

$$
\langle \chi, \alpha^{\vee} \rangle = \frac{2(\chi, \alpha)}{(\alpha, \alpha)}
$$

for all $\chi \in X^*(\mathbb{T})$. We let $\Phi^{\vee} \subseteq X_*(\mathbb{T})$ denote the set of coroots. These are given explicitly as follows. Let $f_i \in X_*(\mathbb{T})$ denote the dual of e_i ; in [\(3.15\)](#page-10-3), it is the coset attached to the *i*-th standard basis vector of \mathbf{Z}^{n+1} . Then $(e_j - e_i)^{\vee}$ = $f_j - f_i$ for $i \neq j$, $(e_0 - e_j - e_i)^{\vee} = -f_j - f_i$ for $i \neq j$, and $(e_0 - 2e_j)^{\vee} = -f_j$. The above gives a description of the root datum $(X^*(\mathbb{T}), \Phi, X_*(\mathbb{T}), \Phi^{\vee})$ of \mathbb{G} = PGSp_{2n}. The complex dual group $\widehat{\mathbb{G}}$ (with dual root datum $(X_*(\mathbb{T}), \Phi^{\vee}, X^*(\mathbb{T}), \Phi))$ is $Spin(2n+1, \mathbb{C})$. In particular, for the maximal torus $\hat{\mathbb{T}}$ of $\hat{\mathbb{G}}$, we have

$$
X^*(\widehat{\mathbb{T}}) \cong X_*(\mathbb{T}), \qquad X_*(\widehat{\mathbb{T}}) \cong X^*(\mathbb{T}). \tag{3.18}
$$

Now define the positive Weyl chamber

$$
\mathcal{C}^+ = \{ \lambda \in X_*(\mathbb{T}) | \langle \chi, \lambda \rangle \ge 0 \text{ for all } \chi \in \Delta \}. \tag{3.19}
$$

We will frequently identify C^+ with its counterpart in $X^*(\hat{\mathbb{T}}) \cong X_*(\mathbb{T})$. By definition, an element $\lambda \in X_*(\mathbb{T})$ belongs to \mathcal{C}^+ if and only if

$$
\langle e_{j+1} - e_j, \lambda \rangle \ge 0
$$
 for $j = 1, ..., n-1$ and $\langle e_0 - 2e_n, \lambda \rangle \ge 0$.

The above holds if and only if, in the notation of [\(3.15\)](#page-10-3), every coset representative $(\ell_0, \ell_1, \ldots, \ell_n)$ for λ satisfies

$$
\ell_1 \le \ell_2 \le \dots \le \ell_n \le \ell_0/2. \tag{3.20}
$$

In fact, each such λ has a unique representative satisfying

$$
0 = \ell_1 \le \ell_2 \le \cdots \le \ell_n \le \ell_0/2. \tag{3.21}
$$

For notational convenience, we will often identify λ with this coset representative.

PROPOSITION 3.1 (Cartan decomposition, [\[Gr\]](#page-67-7)). The group $\mathbb{G}(\mathbf{Q}_p)$ is the disjoint union of the double cosets $K_p \lambda(p) K_p$ for $\lambda \in C^+$. The analogous statement (involving λ satisfying [\(3.20\)](#page-11-0) rather than [\(3.21\)](#page-11-1)) also holds for $G(\mathbf{Q}_p)$.

PROPOSITION 3.2. For $g \in G(\mathbf{Q}_p)$ and $m \geq 1$, let $d_m(g)$ denote the generator (chosen as a power of p) of the fractional ideal of \mathbf{Q}_p generated by the set

 $\{\det B \mid B \text{ is an } m \times m \text{ submatrix of } g\}.$

Then given $\lambda \in X_*(T)$ satisfying [\(3.20\)](#page-11-0), an element $g \in G(\mathbf{Q}_p)$ belongs to the double coset $G(\mathbf{Z}_p)\lambda(p)G(\mathbf{Z}_p)$ if and only if each of the following holds:

- 1. $r(g) = p^{\ell_0}$
- 2. for each $m = 1, ..., n$, $d_m(g) = p^{\ell_1 + ... + \ell_m}$.

Proof. This result can be extracted from Chapter II of Newman [\[Ne\]](#page-69-6). For convenience, we sketch some of the details. Let $g \in G(\mathbf{Q}_p)$. By the Cartan decomposition, there exist $k_1, k_2 \in G(\mathbf{Z}_p)$ such that $k_1 g k_2 = \lambda(p)$ for a unique cocharacter λ satisfying [\(3.20\)](#page-11-0). We need to show that the two conditions given above are satisfied. The converse will then also follow, since $\lambda(p)$ is uniquely determined by its first n diagonal entries and its similitude.

The first condition is immediate. For the second, observe that there exists an integer $a \geq 0$ such that $p^a g \in M_{2n}(\mathbb{Z}_p)$. The *m*-th diagonal coordinate of $k_1p^a g k_2$ is then $p^{a+\ell_m} \in \mathbb{Z}_p$. Regarding $p^a g$ as an element of $M_{2n}(\mathbb{Z}_p)$ and regarding k_1, k_2 as elements of $GL_{2n}(\mathbb{Z}_p)$, by [\[Ne,](#page-69-6) Chap. II, Sect. 16, Eq. (13)] with $R = \mathbf{Z}_p$, we have

$$
p^{a+\ell_m} = \frac{d_m(p^a g)}{d_{m-1}(p^a g)} = \frac{p^{am} d_m(g)}{p^{a(m-1)} d_{m-1}(g)} = p^a \frac{d_m(g)}{d_{m-1}(g)}
$$

(with $d_0(g) = 1$). Hence

$$
p^{\ell_m} = \frac{d_m(g)}{d_{m-1}(g)}.
$$

 \Box

4 Adelic Siegel modular forms

Let **A** denote the adele ring of **Q**, and fix a Haar measure dg on $\mathbb{G}(A)$. Let $L^2 = L^2(\mathbb{G}(\mathbf{Q})\backslash\mathbb{G}(\mathbf{A}))$ be the space of measurable functions $\phi: G(\mathbf{A}) \to \mathbf{C}$ satisfying

• $\phi(z\gamma g) = \phi(g)$ for all $z \in Z(\mathbf{A}), \gamma \in G(\mathbf{Q}), g \in G(\mathbf{A})$

This is easily seen to be equivalent to condition (2), as needed.

$$
\bullet \ \int_{\mathbb{G}({\bf Q})\backslash\mathbb{G}({\bf A})} |\varphi(g)|^2 dg < \infty.
$$

For any parabolic subgroup P of G, P can be written as MN , where M is the Levi subgroup and N is unipotent. An element $\varphi \in L^2$ is cuspidal if for any parabolic subgroup $P = MN$ of G,

$$
\int_{N(\mathbf{Q})\backslash N(\mathbf{A})} \varphi(ng) dn = 0 \quad \text{ for a.e. } g \in G(\mathbf{A}).
$$

We let $L_0^2 \subseteq L^2$ denote the subspace of cuspidal functions.

The right regular representation of $G(\mathbf{A})$ on L^2_0 decomposes discretely as $\bigoplus \pi$ with possible multiplicity allowed, where π are by definition the cuspidal automorphic representations of $G(A)$ with trivial central character. Any such constituent π is a restricted tensor product

$$
\pi = \bigotimes_{p \leq \infty} \pi_p
$$

where π_p is an irreducible admissible representation of $G(\mathbf{Q}_p)$.

Fix an integer $k > n$ with nk even (so that $(A.24)$ is trivial). Then as in $\S A.7$, there is a holomorphic discrete series representation π_k of $\mathbb{G}(\mathbf{R})$ of weight k. In the inner product space π_k , up to unitary scaling there is a unique holomorphic unit vector ϕ_0 satisfying

$$
\pi_{\mathbf{k}}\begin{pmatrix} A & B \\ -B & A \end{pmatrix}\begin{pmatrix} \phi_0 = \det(A + Bi)^{\mathbf{k}} \phi_0 & \begin{pmatrix} A & B \\ -B & A \end{pmatrix} \in K_\infty\end{pmatrix}.\tag{4.1}
$$

Let $\Pi_{\mathbf{k}}(N)$ denote the set of cuspidal representations π of $G(\mathbf{A})$ with trivial central character for which $\pi_{\infty} = \pi_{\mathbf{k}}$ and $\pi_{\text{fin}}^{K_0(N)} \neq 0$. For such π , we let $v_{\pi_{\infty}} \in V_{\pi_{\infty}}$ denote a lowest weight vector as in [\(4.1\)](#page-13-0). For any representation π_{fin} of $G(\mathbf{A}_{fin})$ and any subgroup $U < G(\mathbf{A}_{fin})$, we write π_{fin}^U for the space of U-fixed vectors in the space of π_{fin} . Define

$$
A_{\mathbf{k}}(N) = \bigoplus_{\pi \in \Pi_{\mathbf{k}}(N)} \mathbf{C} v_{\pi_{\infty}} \otimes \pi_{\text{fin}}^{K_0(N)}.
$$
 (4.2)

This corresponds to a classical space of holomorphic Siegel cusp forms of weight k and level N (see §[7.3](#page-27-1) below).

For $\pi \in \Pi_{\mathbf{k}}(N)$, let $E_{\mathbf{k}}(\pi, N)$ be an orthogonal basis for the summand indexed by π in [\(4.2\)](#page-13-1). Then the set

$$
E_{\mathbf{k}}(N) = \bigcup_{\pi \in \Pi_{\mathbf{k}}(N)} E_{\mathbf{k}}(\pi, N) \tag{4.3}
$$

is an orthogonal basis for $A_k(N)$.

5 The test function and its kernel

Any function $f \in L^1(\mathbb{G}(\mathbf{A}))$ defines an operator $R(f)$ on L^2 by

$$
R(f)\phi(x) = \int_{\mathbb{G}(\mathbf{A})} f(g)\phi(xg)dg.
$$

In this section we will define the bi- $K_0(N)$ -invariant test function f to be used in the trace formula.

5.1 Definition of the test function

We define, for $g \in G(\mathbf{R}),$

$$
f_{\infty}(g) = d_{\mathbf{k}} \overline{\langle \pi_{\mathbf{k}}(g)\phi_0, \phi_0 \rangle},
$$

where ϕ_0 is a unit vector in the space of π_k satisfying [\(4.1\)](#page-13-0), and d_k is the formal degree of π_k . The matrix coefficient is independent of the choice of Haar measure on $\mathbb{G}(\mathbf{R})$. It is computed explicitly in Corollary [A.10](#page-61-1) of the Appendix. The formal degree d_k , which depends on the choice of Haar measure, is given in [\(A.21\)](#page-59-1). For our purposes, the particular choice of measure is immaterial. By Proposition [A.6,](#page-58-1) $f_{\infty} \in L^1(\mathbb{G}(\mathbf{R}))$ precisely when

 $k > 2n$.

so this hypothesis will be in force throughout.

We will take $f_{fin} = \prod_{p<\infty} f_p$ to be a bi- $Z(\mathbf{A}_{fin})K_0(N)$ -invariant function on $G(\mathbf{A}_{fin})$, of the following form.^{[†](#page-14-2)} Fix a finite set **S** of prime numbers not dividing N. For $p \in \mathbf{S}$, let f_p be a bi- $Z_p K_p$ -invariant function with compact support modulo Z_p , taking the value 1 on its support. Here, $Z_p = Z(\mathbf{Q}_p)$ is the center of $G(\mathbf{Q}_p)$. By the Cartan decomposition, f_p is the characteristic function of a set of the form Z_pC_p , with C_p a finite union of double cosets of the form $K_p\lambda(p)K_p$, where, by [\(3.21\)](#page-11-1),

$$
\lambda(p) = \text{diag}(p^{\ell_1}, \dots, p^{\ell_n}, p^{\ell_0 - \ell_1}, \dots, p^{\ell_0 - \ell_n})
$$

with $0 = \ell_1 \leq \cdots \leq \ell_n \leq \ell_0/2$. Without any real loss of generality, we make the further assumption that the similitude of C_p has constant valuation, i.e.

$$
r(C_p) = p^{r_p} \mathbf{Z}_p^*
$$

for some integer $r_p \geq 0$. This amounts to requiring that $\ell_0 = r_p$ for each of the $\lambda(p)$ out of which C_p is built. In particular, $C_p = K_p$ if $r_p = 0$. Having fixed

[†]Although we have defined K_p and $K_0(N)_p$ in [\(3.6\)](#page-8-0) as subsets of $\mathbb{G}(\mathbf{Q}_p)$, in this section we will blur the distinction between these sets and their preimages in $G(\mathbf{Z}_p)$. No confusion should occur since everything is invariant under the center.

such f_p for each $p \in \mathbf{S}$, we define the global similitude

$$
r = \prod_{p \in \mathbf{S}} p^{r_p} \ge 1. \tag{5.1}
$$

Now for $p \notin \mathbf{S}$, define $C_p = K_0(N)_p$, and set

$$
\psi(N)_p = \text{meas}(K_0(N)_p)^{-1} = [K_p : K_0(N)_p],
$$

$$
\psi(N) = \prod_p \psi(N)_p = [K_{\text{fin}} : K_0(N)].
$$

We then define $f_p: G(\mathbf{Q}_p) \longrightarrow \mathbf{C}$ by

$$
f_p(g) = \begin{cases} \psi(N)_p & \text{if } g \in Z_p C_p \\ 0 & \text{otherwise.} \end{cases}
$$

(In fact this holds as well when $p \in \mathbf{S}$, since $\psi(N)_p = 1$ in that case.) Having fixed f_{fin} above, we set

$$
f = f_{\infty} \times f_{\text{fin}}.
$$

By our assumption that $k > 2n$, $f \in L^1(\mathbb{G}(A))$. The following is a useful observation about the support of f_{fin} .

PROPOSITION 5.1. Suppose $x, y \in G(\mathbf{A}_{fin})$ satisfy $r(x)^{-1}r(y) \in \mathbf{Z}^*$. Then for $\gamma \in G(\mathbf{Q})$, $f_{fin}(x^{-1}\gamma y) \neq 0$ only if there exists $s \in \mathbf{Q}^*$, uniquely determined up to its sign, such that

$$
r(\gamma) = \pm s^2 r. \tag{5.2}
$$

Suppose $\gamma \in G(\mathbf{Q})$ satisfies [\(5.2\)](#page-15-0), and set $\tilde{\gamma} = s^{-1}\gamma$. Then $r(\tilde{\gamma}) = \pm r$, and $f_{\text{fin}}(x^{-1}\gamma y) \neq 0$ if and only if

$$
x^{-1}\tilde{\gamma}y \in \prod_p C_p = \prod_{p \in \mathbf{S}} C_p \prod_{p|N} K_0(N)_p \prod_{p \nmid N, p \notin \mathbf{S}} K_p \subseteq M_{2n}(\widehat{\mathbf{Z}}). \tag{5.3}
$$

Proof. If $f_{fin}(x^{-1}\gamma y) \neq 0$, then $x^{-1}\gamma y = zc$ for some $z \in Z(\mathbf{A}_{fin})$ and $c \in \mathcal{C}$ $\prod_p C_p$. Since $\mathbf{A}_{\text{fin}}^* = \mathbf{Q}^* \widehat{\mathbf{Z}}^*$, we may write $z = sa$ where $s \in \mathbf{Q}^*$ and $a \in \widehat{\mathbf{Z}}^*$. We may absorb a into c so that $z = s$ without loss of generality. Taking the similitude and using $r(x)^{-1}r(y) \in \mathbb{Z}^*$, we see that $r(\gamma) = s^2ur$ for some $u \in \mathbb{Z}^*$. Since $r(\gamma)$, $s, r \in \mathbf{Q}^*$, it follows that $u \in \mathbf{Z}^* \cap \mathbf{Q}^* = {\pm 1}$, as claimed. It is clear that s is unique up to its sign, and that $x^{-1}\tilde{\gamma}y \in \prod_p C_p$. Conversely, since the support of f_{fin} is $Z(\mathbf{A}_{fin}) \prod_p C_p$, it is also clear that if $x^{-1}\tilde{\gamma}y \in \prod_p C_p$, then $f_{\text{fin}}(x^{-1}\gamma y) \neq 0.$ \Box

5.2 SPECTRAL PROPERTIES OF $R(f)$

Let f be the function defined above. Recall that it is not completely fixed because it depends on some choice of a finite union C_p of double cosets for each prime $p \in S$. Here we show that the Hecke operator $R(f)$ has finite rank, and we compute its effect on adelic Siegel cusp forms.

PROPOSITION 5.2. Let U be a unipotent subgroup of \mathbb{G} and let $g \in \mathbb{G}(\mathbf{R})$. Then for almost all $x \in U(\mathbf{R}) \backslash \mathbb{G}(\mathbf{R}),$

$$
\int_{U(\mathbf{R})} f_{\infty}(gux) du = 0.
$$

Proof. Let U' be a one-dimensional subgroup of U. There exists $E \in M_{2n}(\mathbf{R})$ such that $U'(\mathbf{R}) = \{I + tE \mid t \in \mathbf{R}\}.$ By Corollary [A.10,](#page-61-1) there exist complex numbers A, B and $C \neq 0$ depending on $u \in U(\mathbf{R})$, $g, x \in \mathbb{G}(\mathbf{R})$, and $k > 2n$, such that

$$
f_{\infty}(g(I_n + tE)ux) = \frac{C}{(At+B)^{k}}.
$$

The denominator is nonzero for $t \in \mathbb{R}$ since the matrix coefficient is finite. If $A \neq 0$, then by the fundamental theorem of calculus,

$$
\int_{\mathbf{R}} f_{\infty}(g(I_n + tE)ux)dt = \frac{C}{A}(At + B)^{-k+1}\Big|_{-\infty}^{\infty} = 0.
$$

If $A = 0$, then

$$
\int_{\mathbf{R}} f_{\infty}(g(I_n + tE)ux)dt = \infty.
$$

Hence

$$
\int_{U(\mathbf{R})} f_{\infty}(gux)du = \int_{U'(\mathbf{R}) \backslash U(\mathbf{R})} \int_{U'(\mathbf{R})} f_{\infty}(gu'ux)du'du
$$

is either 0 or divergent. It remains to show that this integral is convergent for almost all $x \in U(\mathbf{R})\backslash \mathbb{G}(\mathbf{R})$. But this is immediate from the fact that because $f_{\infty} \in L^1(\mathbb{G}(\mathbf{R}))$, the integral

$$
\int_{\mathbb{G}(\mathbf{R})} f_{\infty}(gx) dx = \int_{U(\mathbf{R}) \backslash \mathbb{G}(\mathbf{R})} \int_{U(\mathbf{R})} f_{\infty}(gux) du dx
$$

 \Box

is convergent.

PROPOSITION 5.3. For a test function f as defined in $\S 5.1$, and the subspace $A_{\mathbf{k}}(N) \subseteq L_0^2$ given in [\(4.2\)](#page-13-1), $R(f)$ annihilates $A_{\mathbf{k}}(N)^{\perp}$ and maps $A_{\mathbf{k}}(N)$ into itself.

Proof. Complete details for the case of $GL(2)$ are given in [\[KL1,](#page-68-7) Propositions 13.11, 13.12], and everything carries over directly to the case under consideration here, using the above proposition in place of [\[KL1,](#page-68-7) Corollary 13.10]. So

we just briefly sketch the ideas. It follows easily from the above proposition that for any $\phi \in L^2$, $R(f)\phi$ is cuspidal, i.e. $R(f): L^2 \to L_0^2$. Furthermore one checks that f_{∞} is self-dual, so the adjoint $R(f)^*$ also has this property. It then follows that $R(f)$ annihilates $(L_0^2)^{\perp}$. Next, one may use a general orthogonality property of discrete series matrix coefficients, together with the $K_0(N)$ -invariance of f_{fin} , to show that the image of $R(f)$ (and of $R(f)^*$) lies in $A_{k}(N)$. \Box

It remains to compute the effect of $R(f)$ on a nonzero element $v \in A_k(N)$. We may assume without loss of generality that v is a pure tensor in some cuspidal representation $\pi \in \Pi_{\mathbf{k}}(N)$. Write $\pi = \pi_{\mathbf{k}} \otimes \pi' \otimes \bigotimes_{p \in \mathbf{S}} \pi_p$, where π' is a representation of $\prod_{p \notin S} G(Q_p)$. Accordingly, we write $f = f_{\infty} \times f' \times \prod_{p \in S} f_p$ and

$$
v = v_{\infty} \otimes v' \otimes \bigotimes_{p \in \mathbf{S}} v_p,
$$

where $v_{\infty} = \phi_0$ as in [\(4.1\)](#page-13-0). Then (e.g. by [\[KL1,](#page-68-7) Prop. 13.17]) we have

$$
R(f)v = \pi_{\mathbf{k}}(f_{\infty})v_{\infty} \otimes \pi'(f')v' \otimes \bigotimes_{p \in \mathbf{S}} \pi_p(f_p)v_p.
$$

By the orthogonality relations for discrete series matrix coefficients([\[KL1,](#page-68-7) Corollary 10.26]),

$$
\pi_{\mathbf{k}}(f_{\infty})v_{\infty}=v_{\infty},
$$

where $\pi_{\mathbf{k}}(f_{\infty})$ is defined using the same Haar measure as that defining $d_{\mathbf{k}}$. Likewise, because f' is the characteristic function of $K_0(N)' = \prod_{p \notin S} K_0(N)_p$ scaled by meas $(K_0(N)')^{-1}$, and $v' \in \pi'^{K_0(N)'}$, we have

$$
\pi'(f')v' = v'.
$$

For $p \in \mathbf{S}$, since f_p is K_p bi-invariant, $\pi_p(f_p)$ preserves the subspace $\pi_p^{K_p} = \mathbf{C}v_p$. Hence v_p is an eigenvector. Writing $(\mathcal{S}f_p)(t_{\pi_p})$ for the eigenvalue as in [\(2.8\)](#page-7-3), we have

$$
R(f)v = \left(\prod_{p \in \mathbf{S}} (\mathcal{S}f_p)(t_{\pi_p})\right)v
$$
\n(5.4)

for v as above.

5.3 The kernel function

For a test function f as defined in §[5.1,](#page-14-1) the associated kernel function on $G(\mathbf{A}) \times G(\mathbf{A})$ is defined as

$$
K(x, y) = K_f(x, y) = \sum_{\gamma \in \mathbb{G}(\mathbf{Q})} f(x^{-1} \gamma y).
$$

It satisfies

$$
R(f)\phi(x) = \int_{\mathbb{G}(\mathbf{Q})\backslash\mathbb{G}(\mathbf{A})} K(x,y)\phi(y)dy \qquad (\phi \in L^2).
$$

PROPOSITION 5.4. $K(x, y)$ is continuous in both variables. Furthermore, given any subsets $J_1, J_2 \subseteq G(\mathbf{A})$ each having compact image in $\mathbb{G}(\mathbf{A})$, for every $\gamma \in \mathbb{G}(\mathbf{Q})$ there exists a real number α_{γ} , independent of N and k, such that for all $g_1 \in J_1$ and $g_2 \in J_2$,

$$
d_{\mathbf{k}}^{-1}\psi(N)^{-1}|f(g_1^{-1}\gamma g_2)| \le \alpha_\gamma
$$

and

$$
\sum_{\gamma\in\mathbb{G}({\bf Q})}\alpha_\gamma<\infty.
$$

Proof. For $g \in \mathbb{G}(\mathbf{R})$, define $f_{\kappa}(g) = \overline{\langle \pi_{\kappa}(g)\phi_0, \phi_0 \rangle}$, so that $f_{\infty}(g) = d_{\kappa}f_{\kappa}(g)$. Define the function $\tilde{f} = \tilde{f}_{\infty} \tilde{f}_{\text{fin}}$ on $G(\mathbf{A})$, where $\tilde{f}_{\infty} = d_{k_0} f_{k_0}$ with

$$
\mathbf{k}_0 = 2n + 1 \leq \mathbf{k},
$$

and \tilde{f}_{fin} is the characteristic function of $\prod_p Z_p C_p$, where $C_p = K_p$ for all $p \notin S$. In other words, \tilde{f} is the test function we have defined earlier in the special case where $k = k_0$ and $N = 1$. We claim that

$$
\psi(N)^{-1}d_{\mathbf{k}}^{-1}|f(g)| \le d_{\mathbf{k}_0}^{-1}|\tilde{f}(g)|. \tag{5.5}
$$

Since $\text{supp}(f_{fin}) \subseteq \text{supp}(\tilde{f}_{fin}),$ it is clear from the definitions that $|\psi(N)^{-1}|f_{fin}(g)| \leq |\tilde{f}_{fin}(g)| \text{ for all } g \in G(\mathbf{A}_{fin}).$ For the archimedean part, by Corollary [A.10,](#page-61-1) when $g = \begin{pmatrix} A & B \\ C & D \end{pmatrix} \in \mathbb{G}(\mathbf{R})$ we have

$$
\left[\frac{r(g)^{n/2}2^n}{|\det(A+D+i(B-C))|}\right]^k = |f_k(g)| = |\langle \pi_k(g)\phi_0, \phi_0 \rangle|
$$

$$
\leq ||\pi_k(g)\phi_0|| ||\phi_0|| = 1,
$$

where we have used the fact that ϕ_0 is a unit vector and π_k is unitary. It follows that the expression in the brackets is at most 1, and hence $|f_{\mathbf{k}}(g)| \leq |f_{\mathbf{k}_0}(g)|$ for all $q \in \mathbb{G}(\mathbf{R})$. This proves [\(5.5\)](#page-18-0).

Hence, it suffices to prove the assertion for $f = \tilde{f}$ (with $N = 1$). By [\[Li,](#page-69-3) Prop. 3.1], it suffices to show that there exist bounded compactly supported functions ψ_1, ψ_2 on $\mathbb{G}(\mathbf{A})$ such that $\tilde{f} = \tilde{f} * \psi_1 = \psi_2 * \tilde{f}$, where convolution is defined by

$$
f_1 * f_2 = \int_{\mathbb{G}(\mathbf{A})} f_1(\alpha) f_2(\alpha^{-1} g) d\alpha = \int_{\mathbb{G}(\mathbf{A})} f_1(g \alpha^{-1}) f_2(\alpha) d\alpha.
$$

For the finite part of ψ_i , we take the characteristic function of $Z(\mathbf{A}_{fin})K_0(N)$, multiplied by the reciprocal of the measure of this set in $\mathbb{G}(\mathbf{A}_{fin})$. It remains to

define the archimedean components of ψ_1, ψ_2 . We claim first that there exists a measurable compactly supported function ξ on $\mathbb{G}(\mathbf{R})$ satisfying

$$
\xi(\begin{pmatrix} U & V \\ -V & U \end{pmatrix} g) = \det(U + iV)^{k_0} \xi(g)
$$

for all $\begin{pmatrix} U & V \\ -V & U \end{pmatrix} \in K_{\infty}$, or equivalently, since $G(\mathbf{R}) = K_{\infty} P(\mathbf{R})$,

$$
\xi(\begin{pmatrix}U&V\\-V&U\end{pmatrix}\begin{pmatrix}A&B\\O&D\end{pmatrix})=\det(U+iV)^{\mathtt{k}_0}\xi(\begin{pmatrix}A&B\\O&D\end{pmatrix}).
$$

Let $P' \subseteq P(\mathbf{R})$ be a set of representatives for $(K_{\infty} \cap P(\mathbf{R})) \backslash P(\mathbf{R})$. Then every element $g \in G(\mathbf{R})$ has a unique decomposition $g = kb'$ with $k \in K_{\infty}$ and $b' \in P'$. We may choose P' so that it is a real manifold, and then take ξ to be any compactly supported function on P', extended to $G(\mathbf{R})$ by $\xi(kb') =$ $\det(k)^{k_0}\xi(b')$. One choice of P' is given as follows. Notice that

$$
K_{\infty} \cap P(\mathbf{R}) = \{ \begin{pmatrix} U & \\ & U \end{pmatrix} | U^t U = I \} \cong \mathcal{O}(n),
$$

so using the decomposition (3.5) of $P(\mathbf{R})$, we see that

$$
(K_{\infty} \cap P(\mathbf{R})) \backslash P(\mathbf{R}) \cong \left(\mathrm{O}(n) \backslash \mathrm{GL}_n(\mathbf{R})\right) \times \mathbf{R}^* \times S_n(\mathbf{R}).
$$

Hence we can take $P' \subseteq P(\mathbf{R})$ to be the subgroup identifying as

$$
P' \cong \Big\{ \left(\begin{array}{cccc} a_{11} & * & \dots & * \\ & & \ddots & \\ & & \ddots & \\ & & & a_{nn} \end{array} \right) \mid a_{jj} > 0 \Big\} \times \mathbf{R}^* \times S_n(\mathbf{R}).
$$

The proof now proceeds exactly as in [\[Li,](#page-69-3) Prop. 3.2]: one sees easily that $\pi_{k_0}(\xi)\phi_0$ is a vector of weight k_0 as in [\(4.1\)](#page-13-0), and hence must be a multiple of ϕ_0 . For an appropriate choice of ξ , the multiple is nonzero, and hence without loss of generality $\pi_{k_0}(\xi)\phi_0 = \phi_0$. From here it is easy to show directly that $\xi * f_{k_0} = f_{k_0}$, so we can take $\psi_1 = \xi \times \psi_{fin}$. Similarly, we can take $\psi_2 = \psi_1^*$.

5.4 Spectral expression for the kernel

Because the support of the test function is not compact modulo the center, some care is needed in order to justify the spectral expansion of the kernel function.

PROPOSITION 5.5. With notation as in \S_4 , the kernel function of the operator $R(f)$ has the spectral expansion

$$
K(x,y) = \sum_{\pi \in \Pi_{\mathbf{k}}(N)} \left(\sum_{\varphi \in E_{\mathbf{k}}(\pi,N)} \frac{\varphi(x)\overline{\varphi(y)}}{\|\varphi\|^2} \right) \prod_{p \in \mathbf{S}} (\mathcal{S}f_p)(t_{\pi_p}). \tag{5.6}
$$

Proof. As shown in §[5.2,](#page-16-0) the operator $R(f)$ vanishes on $A_k(N)^{\perp}$ and is diagonalizable on $A_k(N)$, the elements of $E_k(N)$ (see [\(4.3\)](#page-13-2)) being eigenvectors. It follows easily that

$$
\Phi(x,y) = \sum_{\varphi \in E_{\mathbf{k}}(N)} \frac{R(f)\varphi(x)\varphi(y)}{\|\varphi\|^2}
$$

is a kernel function for $R(f)$. (See e.g. [\[KL1,](#page-68-7) p. 228] for details.) It follows that $K(x, y) = \Phi(x, y)$ a.e. On the other hand, Φ is continuous in both variables, being a sum of finitely many adelic Siegel cusp forms, while by Proposition [5.4,](#page-18-1) $K(x, y)$ is also continuous. Hence they are equal everywhere. Using [\(5.4\)](#page-17-1) we see that $\Phi(x, y)$ is equal to the spectral expression given in [\(5.6\)](#page-19-1). \Box

6 ASYMPTOTIC FOURIER TRACE FORMULA FOR $GSp(2n)$

6.1 Additive characters

Let

$$
\theta: \mathbf{Q} \backslash \mathbf{A} \longrightarrow \mathbf{C}^*
$$

be the nontrivial character whose local components are given by

$$
\theta_p(x) = \begin{cases} e^{2\pi ix} & \text{if } p = \infty \\ e^{-2\pi ir_p(x)} & \text{if } p < \infty, \end{cases}
$$
\n(6.1)

where $r_p(x) \in \mathbf{Q}$ is any number with p-power denominator satisfying $x \in$ $r_p(x) + \mathbf{Z}_p$. All characters of the additive group $\mathbf{Q} \setminus \mathbf{A}$ are of the form $\theta(qx)$ for some $q \in \mathbf{Q}$. It follows easily that any character of the additive group $S_n(\mathbf{Q})\backslash S_n(\mathbf{A})$ has the form $S \mapsto \theta(\text{tr }\sigma S)$ for some $\sigma \in S_n(\mathbf{Q})$. We fix two such matrices $\sigma_1, \sigma_2 \in S_n(\mathbf{Q})$ and define

$$
\theta_j(S) = \theta(\text{tr }\sigma_j S) \qquad (j = 1, 2).
$$

Since we are interested in the θ_j -Fourier coefficients of Siegel cusp forms, we can in fact assume that the σ_j belong to the set

 \mathcal{R}_n^+ = set of half-integral positive definite symmetric matrices $\sigma \in GL_n(\mathbf{Q})$

(this means 2σ has integer entries and even diagonal entries).

6.2 The setup and the spectral side

Given a continuous function φ on $\mathbb{G}(Q)\backslash \mathbb{G}(A)$, its Whittaker function along the unipotent subgroup $N = \{$ $\int I_n \hat{S}$ O I_n Ñ $| S \in S_n(A) \}$ is defined by

$$
W_{\varphi}(g,\chi) = \int_{N(\mathbf{Q}) \backslash N(\mathbf{A})} \varphi(ng) \overline{\chi(n)} dn
$$

for $g \in \mathbb{G}(\mathbf{A})$ and χ a character. Write $n_S = \begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix}$ for $S \in S_n(\mathbf{A})$. There exits $\sigma \in S_n(\mathbf{Q})$ such that $\chi(n_S) = \theta(\text{tr }\sigma S)$ for all S. Define

$$
c_{\sigma}(\varphi) = W(1, \chi) = \int_{S_n(\mathbf{Q}) \backslash S_n(\mathbf{A})} \varphi(n_S) \overline{\theta(\text{tr }\sigma S)} dS. \tag{6.2}
$$

In §[5.1,](#page-14-1) we defined a test function, henceforth to be denoted f_N , from the following data: a finite set S of primes, a level N coprime to S , a compact set $C_p \subseteq G(\mathbf{Q}_p)$ for each $p \in \mathbf{S}$, and a weight $k > 2n$. In this section, we compute the following limit:

$$
I = \lim_{\substack{N \to \infty \\ (N, \mathbf{S}) = 1}} \iint_{(N(\mathbf{Q}) \backslash N(\mathbf{A}))^2} \frac{K_{f_N}(n_1, n_2)}{\psi(N)} \overline{\theta_1(n_1)} \theta_2(n_2) dn_1 dn_2, \tag{6.3}
$$

where $\psi(N) = \text{meas}(\overline{K_0(N)})^{-1}$. We will compute I in two ways, using the spectral and geometric forms of the kernel function. The preliminary form of the resulting formula is given in Theorem [6.3](#page-24-1) below.

Using the spectral form [\(5.6\)](#page-19-1) of the kernel function along with [\(6.2\)](#page-21-1), we formally obtain

$$
I = \lim_{\substack{N \to \infty \\ (N, \mathbf{S}) = 1}} \frac{1}{\psi(N)} \sum_{\pi \in \Pi_{\mathbf{k}}(N)} \sum_{\varphi \in E_{\mathbf{k}}(\pi, N)} \frac{c_{\sigma_1}(\varphi) \overline{c_{\sigma_2}(\varphi)}}{\|\varphi\|^2} \prod_{p \in S} (\mathcal{S}f_p)(t_{\pi_p}). \tag{6.4}
$$

The existence of the limit will be demonstrated in Proposition [6.1](#page-21-2) below.

6.3 THE GEOMETRIC SIDE

First we show how the limit can be eliminated on the geometric side.

PROPOSITION 6.1. Let $f = f_1$ be the test function we have defined when $N = 1$. (Its finite part is the characteristic function of $\prod Z_p C_p$, where $C_p = K_p$ for all $p \notin S$.) Then for I as in [\(6.3\)](#page-21-3), the limit exists, and

$$
I = \iint_{(N(\mathbf{Q})\backslash N(\mathbf{A}))^2} \sum_{\gamma \in \mathbb{P}(\mathbf{Q})} f(n_1^{-1}\gamma n_2) \overline{\theta_1(n_1)} \theta_2(n_2) dn_1 dn_2.
$$
 (6.5)

Proof. Recall that $[0, 1) \times \hat{\mathbb{Z}}$ is a fundamental domain in **A** for $\mathbb{Q}\backslash\mathbf{A}$. We may therefore replace $N(\mathbf{Q})\backslash N(\mathbf{A})\cong N(\mathbf{Q}\backslash \mathbf{A})$ by the compact set $J=N([0,1]\times \widehat{\mathbf{Z}})$. Applying Proposition [5.4](#page-18-1) with $J_1 = J_2 = J$, the integrand in [\(6.3\)](#page-21-3) is absolutely bounded by the constant $\sum \alpha_{\gamma}$. Hence by the dominated convergence theorem,

$$
I = \lim_{N \to \infty \atop (N,S)=1} \iint_{J \times J} \frac{K_{f_N}(n_1, n_2)}{\psi(N)} \overline{\theta_1(n_1)} \theta_2(n_2) dn_1 dn_2
$$

$$
=\iint_{J\times J} \sum_{\gamma\in \mathbb{G}({\bf Q})}\lim_{(N,\infty)\atop{(N,S)=1}}\frac{f_N\big(n_1^{-1}\gamma n_2\big)}{\psi(N)}\overline{\theta_1(n_1)}\theta_2(n_2)dn_1dn_2.
$$

By Proposition [5.1,](#page-15-1) we can assume $r(\gamma) = \pm r$. Furthermore, if

$$
f_{N,\text{fin}}(n_{1\,\text{fin}}^{-1}\gamma n_{2\,\text{fin}}) \neq 0
$$

for $n_{i \text{ fin}} \in N(\hat{\mathbf{Z}})$, then

$$
n_{1\,\text{fin}}^{-1}\gamma n_{2\,\text{fin}}\in\prod_{p\in\mathbf{S}}C_p\prod_{p\notin\mathbf{S}}K_0(N)_p.
$$

Therefore

$$
\gamma = \begin{pmatrix} W & X \\ Y & Z \end{pmatrix} \in \prod_{p \in \mathbf{S}} C_p \prod_{p \notin \mathbf{S}} K_0(N)_p \subseteq M_{2n}(\widehat{\mathbf{Z}}).
$$

In particular $Y \in M_n(N\mathbf{Z})$. Hence for any $Y \neq 0$, $f_{N,\text{fin}}(n_{1\text{fin}}^{-1}\gamma n_{2\text{fin}}) = 0$ when N is sufficiently large. On the other hand, if $Y = 0$, then for $n_j \in J$,

$$
\frac{f_N(n_1^{-1}\gamma n_2)}{\psi(N)} = f(n_1^{-1}\gamma n_2) = f_{\infty}(n_{1,\infty}^{-1}\gamma n_{2,\infty})f_{\text{fin}}(\gamma)
$$

for $f = f_1$, which is obviously independent of N. In particular,

$$
I = \iint_{J \times J} \sum_{\gamma \in \mathbb{P}(\mathbf{Q})} f_1(n_1^{-1} \gamma n_2) \overline{\theta_1(n_1)} \theta_2(n_2) dn_1 dn_2.
$$

As a function of n_1 and n_2 , the summation over $\mathbb{P}(\mathbf{Q})$ is $N(\mathbf{Q})$ -invariant in both variables. So we can replace the region of the double integral by $(N(Q)\backslash N(A))^2$. This completes the proof. \Box

Remark: As indicated in the Introduction, the same idea cannot directly be used to understand the case of fixed N and $k \to \infty$. The above argument hinges on the fact that

$$
\lim_{N \to \infty} \frac{f(n_1^{-1} \gamma n_2)}{\psi(N)} = 0
$$

for $n_1, n_2 \in J$ and $\gamma \notin \mathbb{P}(\mathbf{Q})$. For the case of varying k with $N = 1$, the analogous argument would require

$$
\lim_{k \to \infty} \frac{f(n_1^{-1} \gamma n_2)}{c_{nk\sigma}} = 0 \tag{6.6}
$$

for the same γ , n_1 , n_2 , with $c_{n\mathbf{k}\sigma}$ as in Theorem [1.1.](#page-3-0) However, this is false. We have $c_{n\mathbf{k}\sigma}^{-1} =$ $CB^{\mathbf{k}}\prod_{j=1}^n\Gamma(\mathbf{k}-\frac{n+j}{2})$ $\frac{d_k}{d_k}$ for positive constants B and C depending only on n, and d_k is a polynomial in k. By Stirling's approximation, the

product of Gamma factors grows faster than $\mathbf{k}^{(n-\varepsilon)\mathbf{k}}$ for any $\varepsilon > 0$. On the other hand, as we will show in Corollary [A.11,](#page-62-0) $d_{\mathbf{k}}^{-1}|f(n_1\gamma n_2)|$ has at most exponential decay as $k \to \infty$. It follows that the limit appearing in [\(6.6\)](#page-22-0) is actually divergent for many choices of $\gamma \in \mathbb{G}(\mathbf{Q}) \setminus \mathbb{P}(\mathbf{Q})$.

We now examine the sum in [\(6.5\)](#page-21-4). By Proposition [5.1,](#page-15-1) we may assume that $r(\gamma) = \pm r$. Because $f_{\infty}(g) = 0$ if $r(g) < 0$, we may in fact take $r(\gamma) = r$. By a variant of (3.5) , we may write

$$
\gamma = \eta \, g_{A,r},
$$

where $\eta \in N(\mathbf{Q})$ and, for $A \in GL_n(\mathbf{Q}),$

$$
g_{A,r}:=\begin{pmatrix} A&\\&r\,{}^tA^{-1}\end{pmatrix}.
$$

Recalling that $\gamma \in \mathbb{P}(\mathbf{Q})$ is only defined up to multiplication by the center, we observe that for $\lambda \in \mathbf{Q}^*$, $\lambda g_{A,r} = g_{\lambda A,r}$. Hence by our insistance that $r(\gamma) = r$, we may only scale by $\lambda = \pm 1$. Therefore

$$
\sum_{\gamma \in \mathbb{P}(\mathbf{Q})} f(n_1^{-1}\gamma n_2) = \sum_{A \in \mathrm{GL}_n(\mathbf{Q})/\{\pm I_n\}} \sum_{\eta \in N(\mathbf{Q})} f(n_1^{-1}\eta g_{A,\tau} n_2).
$$

Taking, as we may, n_1, n_2 to range through the fundamental domain $N([0, 1] \times$ $\hat{\mathbf{Z}}$) for $N(\mathbf{Q})\backslash N(\mathbf{A})$, by [\(5.3\)](#page-15-2) a given summand vanishes unless

$$
\eta g_{A,r} \in n_{1\,\text{fin}} M_{2n}(\widehat{\mathbf{Z}}) n_{2\,\text{fin}}^{-1} \subseteq M_{2n}(\widehat{\mathbf{Z}}).
$$

This implies that $A \in M_n(\mathbf{Z})$ and $rA^{-1} \in M_n(\mathbf{Z})$. Hence

$$
I = \sum_{\substack{A \in \{\pm I_n\} \setminus M_n(\mathbf{Z})\\rA^{-1} \in M_n(\mathbf{Z})}} \iint_{(S_n(\mathbf{Q}) \setminus S_n(\mathbf{A}))^2} \sum_{S \in S_n(\mathbf{Q})} f\left(\begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix} \begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix} g_{A,r} \begin{pmatrix} I_n & S_2 \\ O & I_n \end{pmatrix}\right)
$$

$$
\times \theta (-\operatorname{tr} \sigma_1 S_1 + \operatorname{tr} \sigma_2 S_2) dS_1 dS_2.
$$

The double integral becomes

$$
\int_{S_n(\mathbf{Q})\backslash S_n(\mathbf{A})} \int_{S_n(\mathbf{A})} f\left(\left(\begin{smallmatrix} I_n & -S_1 \\ O & I_n \end{smallmatrix}\right) g_{A,r}\left(\begin{smallmatrix} I_n & S_2 \\ O & I_n \end{smallmatrix}\right)\right) \theta(-\operatorname{tr} \sigma_1 S_1 + \operatorname{tr} \sigma_2 S_2) dS_1 dS_2,
$$

which equals

$$
\int_{S_n(\mathbf{Q})\backslash S_n(\mathbf{A})} \int_{S_n(\mathbf{A})} f\left(\left(\begin{smallmatrix} I_n & -(S_1 - r^{-1}AS_2 \,^t A) \\ 0 & I_n \end{smallmatrix} \right) g_{A,r} \right) \theta(-\operatorname{tr} \sigma_1 S_1 + \operatorname{tr} \sigma_2 S_2) dS_1 dS_2.
$$

Making the substitution $S_1' = S_1 - r^{-1} A S_2^{\dagger} A$, the above becomes

$$
\iint_{S_n(\mathbf{A}) \times (S_n(\mathbf{Q}) \backslash S_n(\mathbf{A}))} f(\left(\begin{pmatrix} I_n & -S'_1 \\ O & I_n \end{pmatrix} g_{A,r} \right)
$$

\$\times \theta(-\text{tr}(\sigma_1 S'_1 + r^{-1} \sigma_1 A S_2 \, tA) + \text{tr} \, \sigma_2 S_2) dS'_1 dS_2\$
=
$$
\int_{S_n(\mathbf{A})} f(\left(\begin{pmatrix} I_n & -S'_1 \\ O & I_n \end{pmatrix} g_{A,r} \right) \theta(-\text{tr} \, \sigma_1 S'_1) dS'_1
$$

\$\times \int_{S_n(\mathbf{Q}) \backslash S_n(\mathbf{A})} \theta(\text{tr}(\sigma_2 S_2 - r^{-1} \sigma_1 A S_2 \, tA)) dS_2\$.

The value of the second integral is

$$
\begin{cases}\n\text{meas}(S_n(\mathbf{Q}) \setminus S_n(\mathbf{A})) = 1 & \text{if } \theta(\text{tr}(\sigma_2 S_2 - r^{-1} \sigma_1 A S_2 \,^t A)) = 1 \\
\text{for all } S_2 \in S_n(\mathbf{A}), \\
0 & \text{otherwise.} \n\end{cases}
$$

LEMMA 6.2. We have $\theta(\text{tr}(\sigma_2S - r^{-1}\sigma_1AS^tA)) = 1$ for all $S \in S_n(\mathbf{A})$ if and only if

$$
{}^t\!A\sigma_1A = r\sigma_2.
$$

Proof. By the fact that $tr(AB) = tr(BA)$,

$$
\theta(\text{tr}(\sigma_2 S - r^{-1}\sigma_1 AS^t A)) = \theta(\text{tr}((\sigma_2 - r^{-1}{}^t A \sigma_1 A)S)).
$$

The lemma follows from this.

 \Box

6.4 Main formula

The result of the above computation is the following asymptotic Petersson formula for $PGSp(2n)$.

THEOREM 6.3. Let **S** be a finite set of prime numbers, fix sets C_p as in §[5.1,](#page-14-1) and let $r \geq 1$ be the integer defined in [\(5.1\)](#page-15-3). Let $f = f_1$ be the associated test function defined there for the case of level $N = 1$. Then for $k > 2n$ with $2|nk$, and any $\sigma_1, \sigma_2 \in \mathcal{R}_n^+$,

$$
\lim_{\substack{N \to \infty \\ (N,S)=1}} \frac{1}{\psi(N)} \sum_{\pi \in \Pi_k(N)} \left(\sum_{\varphi \in E_k(\pi, N)} \frac{c_{\sigma_1}(\varphi) \overline{c_{\sigma_2}(\varphi)}}{\|\varphi\|^2} \right) \prod_{p \in \mathbf{S}} (\mathcal{S}f_p)(t_{\pi_p})
$$
\n
$$
= \sum_A \int_{S_n(\mathbf{A})} f\left(\begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix} \begin{pmatrix} A & O \\ O & r^t A^{-1} \end{pmatrix} \right) \theta(\text{tr } \sigma_1 S) dS. \tag{6.7}
$$

Here, $\Pi_{\mathbf{k}}(N)$ and $E_{\mathbf{k}}(\pi, N)$ were defined in \S_4 , $S_n(\mathbf{A})$ is the set of symmetric $n \times n$ matrices over the adeles **A**, and A runs through the finite set of rank n matrices in $M_n(\mathbf{Z})/\{\pm I_n\}$ satisfying the following conditions:

1. $r^t A^{-1} \in M_n(\mathbf{Z}).$

2. ${}^tA\sigma_1A = r\sigma_2$.

Remarks: (1) In Appendix [B,](#page-63-0) we give a quantitative version of the above for $GSp(4)$.

(2) Haar measure is normalized as follows. On the geometric side and in the def-inition of the Fourier coefficients [\(6.2\)](#page-21-1) we take meas($S_n(\hat{\mathbf{Z}})$) = 1 and Lebesgue measure on the Euclidean space $S_n(\mathbf{R})$. The Satake transform is defined using meas(K_p) = 1. The archimedean test function f_{∞} depends (via $d_{\bf k}$) on an unspecified choice of measure on $\mathbb{G}(\mathbf{R})$. This choice materializes on the spectral side in $\|\varphi\|^{-2}$. The exact relationship between several natural choices of Haar measure on $Sp_{2n}(\mathbf{R})$ is computed in [\[PSS,](#page-69-7) §A].

Proof. In view of the discussion in the previous two subsections, it just remains to prove that only finitely many matrices A satisfy the given conditions. For fixed j, let $d \in \mathbf{R}$ be the entry in the j-th row and j-th column of $r\sigma_2$. Suppose $A \in M_n(\mathbf{Z})$ satisfies condition 2, and let $v \in \mathbf{Z}^n$ denote the j-th column of A. Then $t_{v\sigma_1v} = d$. We will show that there are only finitely many such v (this is well-known), from which it follows that the set of A is also finite since i is arbitrary.

Because σ_1 is symmetric, there exists an orthogonal matrix Q such that σ_1 = ^tQΛQ, where $\Lambda = \text{diag}(\lambda_1, \ldots, \lambda_n)$ is diagonal. Furthermore the eigenvalues λ_i are all positive since σ_1 is positive definite. It follows that the linear map $Q: \mathbf{R}^n \longrightarrow \mathbf{R}^n$ restricts to give an isometry between the sets

$$
X = \{ v \in \mathbf{R}^n \mid {}^t v \sigma_1 v = d \}
$$

and

$$
Y = \{ w \in \mathbf{R}^n \mid {}^t w \Lambda w = d \}.
$$

Notice that Y is the ellipsoid

$$
\lambda_1 x_1^2 + \dots + \lambda_n x_n^2 = d,
$$

which is compact. Hence X is also compact. Since \mathbf{Z}^n is discrete, it follows that there are only finitely many integer lattice points in X , as claimed. \Box

7 Refinement of the geometric side

Fix a matrix $A \in M_n(\mathbf{Z})$ satisfying the hypotheses of Theorem [6.3.](#page-24-1) The integral on the geometric side of [\(6.7\)](#page-24-2) can be factorized as

$$
I_A = \prod_{p \le \infty} I_{A,p} = \prod_{p \le \infty} \int_{S_n(\mathbf{Q}_p)} f_p(\begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix} \begin{pmatrix} A & O \\ O & r^t A^{-1} \end{pmatrix}) \theta_p(\text{tr } \sigma_1 S) dS.
$$

We will see that for all $p \notin S$, $I_{A,p}$ can be computed explicitly, and is independent of A when k is even.

7.1 Archimedean integral

When $p = \infty$, by Corollary [A.10](#page-61-1) in the Appendix,

$$
I_{A,\infty} = \int_{S_n(\mathbf{R})} f_{\infty}(\begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix} \begin{pmatrix} A & O \\ O & r^t A^{-1} \end{pmatrix}) e^{2\pi i \operatorname{tr} \sigma_1 S} dS
$$
\n
$$
= \int_{S_n(\mathbf{R})} \frac{d_{\mathbf{k}} 2^{n\mathbf{k}} r^{n\mathbf{k}/2} e^{2\pi i \operatorname{tr} \sigma_1 S}}{\det(A + r^t A^{-1} + i S r^t A^{-1})^{\mathbf{k}}} dS
$$
\n
$$
= \frac{d_{\mathbf{k}} 2^{n\mathbf{k}} (\det A)^{\mathbf{k}}}{r^{n\mathbf{k}/2}} \int_{S_n(\mathbf{R})} \frac{e^{2\pi i \operatorname{tr} \sigma_1 S}}{\det(r^{-1} A^t A + I_n + i S)^{\mathbf{k}}} dS. \tag{7.1}
$$

We apply the following formula of Ingham (also found in a paper of Siegel). Let

$$
\Gamma_n(a) = \pi^{n(n-1)/4} \prod_{j=1}^n \Gamma(a - \frac{1}{2}(j-1)).
$$
\n(7.2)

Then for $\delta > \frac{n-1}{2}$ and symmetric matrices $X_0, \Lambda > 0$,

$$
\int_{S_n(\mathbf{R})} \frac{e^{\text{tr} i\Lambda Y}}{\det(X_0 + iY)^{\delta + (n+1)/2}} dY
$$
\n
$$
= \frac{1}{i^{n(n+1)/2} 2^{(n-1)n/2}} \frac{(2\pi i)^{n(n+1)/2} (\det \Lambda)^{\delta}}{\Gamma_n(\delta + (n+1)/2)} e^{-\text{tr}\Lambda X_0}
$$
\n
$$
= \frac{2^n \pi^{n(n+1)/2} (\det \Lambda)^{\delta}}{\Gamma_n(\delta + (n+1)/2)} e^{-\text{tr}\Lambda X_0}
$$
\n(7.3)

(cf. $[I, (1)]$, $[Si1, Hilfssatz 37, p. 585]$; we have used the form given by Herz $[H,$ (1.2)]). In [\[H\]](#page-67-9), the measure is $dZ = \prod_{j \leq k} dz_{jk}$, where $Z = X_0 + iY = (\eta_{jk}z_{jk})$ with $\eta_{jk} = 1$ if $j = k$, and $1/2$ otherwise. Thus

$$
dZ = i^{n(n+1)/2} 2^{(n-1)n/2} dY,
$$

which explains the first factor in the above formula.

We evaluate [\(7.1\)](#page-26-1) using [\(7.3\)](#page-26-2) with $\delta = k - (n + 1)/2 > (n - 1)/2$, $X_0 =$ $I_n + r^{-1}A^tA$, and $\Lambda = 2\pi\sigma_1$. By condition (2) of Theorem [6.3,](#page-24-1) $det(A)$ $\pm r^{n/2} \left(\frac{\det \sigma_2}{\det \sigma_1} \right)$ $\int^{1/2}$. So [\(7.1\)](#page-26-1) becomes

$$
\frac{d_{\bf k} 2^{n{\bf k}}}{r^{nk/2}} \, {\rm sgn}(\det A)^{\bf k} r^{nk/2} \left(\frac{\det \sigma_2}{\det \sigma_1} \right)^{{\bf k}/2} 2^n \pi^{n(n+1)/2} \frac{\det (2\pi \sigma_1)^{{\bf k}-(n+1)/2}}{\Gamma_n({\bf k}) \, e^{2\pi \, \text{tr}(\sigma_1(I_n+r^{-1}A^tA))}}.
$$

We simplify the above using $\text{tr}(\sigma_1 r^{-1} A^t A) = \text{tr}(r^{-1}{}^t A \sigma_1 A) = \text{tr} \sigma_2$ for matrices A as in Theorem [6.3.](#page-24-1) The result is the following:

$$
I_{A,\infty} = \text{sgn}(\det A)^k d_k \left(\frac{\det \sigma_2}{\det \sigma_1} \right)^{k/2} \frac{(\det \sigma_1)^{k-(n+1)/2} (4\pi)^{nk}}{2^{n(n-1)/2} \Gamma_n(\mathbf{k}) e^{2\pi \text{ tr}(\sigma_1 + \sigma_2)}}.
$$
(7.4)

Observe that this is independent of A when k is even.

7.2 NONARCHIMEDEAN INTEGRALS: $p \notin \mathbf{S}$

In this case, $r \in \mathbb{Z}_p^*$, and $A \in GL_n(\mathbb{Z}_p)$. It follows that

$$
f_p(\begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix} \begin{pmatrix} A & O \\ O & r^t A^{-1} \end{pmatrix}) \neq 0 \iff \begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix} \begin{pmatrix} A & O \\ O & r^t A^{-1} \end{pmatrix} \in K_p
$$

$$
\iff S \in M_n(\mathbf{Z}_p).
$$

For such S, $\theta_p(\text{tr }\sigma_1S) = 1$, so we find that

$$
I_{A,p} = \int_{S_n(\mathbf{Q}_p)} f_p(\begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix} \begin{pmatrix} A & O \\ O & r^t A^{-1} \end{pmatrix}) \theta_p(\text{tr }\sigma_1 S) dS = \text{meas}(S_n(\mathbf{Z}_p)) = 1.
$$

(Recall that we use the $N = 1$ test function f_1 in [\(6.7\)](#page-24-2), so that f_p is the characteristic function of $Z_p K_p$ when $p \notin \mathbf{S}$.)

7.3 Proof of Theorem [1.1](#page-3-0)

When $S = \emptyset$, we now have a completely explicit expression for the right-hand side of [\(6.7\)](#page-24-2). Since $r = 1$, the sum over A is nonzero only if ${}^t\!A\sigma_1A = \sigma_2$ for some $A \in GL_n(\mathbf{Z})/\{\pm I_n\}$. In particular, $(\det A)^2 = 1$, so $\det \sigma_1 = \det \sigma_2$. If we let I_A denote the summand indexed by A in Theorem [6.3,](#page-24-1) then by the above discussion,

$$
I_A = (\det A)^k d_k (4\pi)^{nk - n(n-1)/4} \frac{(\det \sigma_1)^{k - (n+1)/2}}{\prod_{j=1}^n \Gamma(k - \frac{j-1}{2})} e^{-2\pi \operatorname{tr}(\sigma_1 + \sigma_2)}.
$$
 (7.5)

We wish to express the spectral side in classical terms. With notation as in (3.3) , let $S_k(N)$ be the space of Siegel cusp forms F satisfying

$$
F(\gamma \cdot \mathfrak{Z}) = j(\gamma, Z)^{k} F(\mathfrak{Z})
$$

for all $\mathfrak{Z} \in \mathcal{H}_n$ and

$$
\gamma \in \Gamma_0(N) = \{ \left(\begin{smallmatrix} A & B \\ C & D \end{smallmatrix} \right) \in \text{Sp}_{2n}(\mathbf{Z}) | C \equiv O \mod N \},\
$$

where

$$
j(g, 3) = r(g)^{-n/2} \det(C_3 + D)
$$
 $(g = \begin{pmatrix} A & B \\ C & D \end{pmatrix} \in G(\mathbf{R})$).

Any $F \in S_k(N)$ has a Fourier expansion

$$
F(3) = \sum_{\sigma \in \mathcal{R}_n^+} a_{\sigma}(F) e^{2\pi i \operatorname{tr} \sigma_3}, \quad (3 \in \mathcal{H}_n). \tag{7.6}
$$

We normalize the Petersson/Maass scalar product on $S_k(N)$ by

$$
\langle F, H \rangle = \frac{1}{\psi(N)} \int_{\Gamma_0(N) \backslash \mathcal{H}_n} F(3) \overline{H(3)} (\det Y)^{k-n-1} dX dY \quad (3 = X + iY), \tag{7.7}
$$

where $\psi(N) = [K_{\text{fin}} : K_0(N)] = [\Gamma_0(1) : \Gamma_0(N)].$

We need to choose the quotient measure on $\mathbb{G}(Q)\backslash\mathbb{G}(A)$ compatibly with the above. For any $N \geq 1$, let $D_N \subseteq \mathcal{H}_n$ be a fundamental domain for $\Gamma_0(N) \backslash \mathcal{H}_n$, identified with a subset of $Sp_{2n}(\mathbf{R})$ via

$$
\mathfrak{Z} = X + iY \longleftrightarrow b_{\mathfrak{Z}} = \begin{pmatrix} Y^{1/2} & XY^{-1/2} \\ O & Y^{-1/2} \end{pmatrix}.
$$

Then, as in [\[KL1,](#page-68-7) Prop. 7.43] for example, we may define a quotient measure on $\mathbb{G}(Q)\backslash\mathbb{G}(\mathbf{A})$ by

$$
\int_{\mathbb{G}(\mathbf{Q})\backslash\mathbb{G}(\mathbf{A})} h(g) dg = \int_{D_N K_{\infty} \times K_0(N)} h(b_3 k_{\infty} \times k_{\text{fin}}) \frac{dX dY}{(\det Y)^{n+1}} dk_{\infty} dk_{\text{fin}},
$$

where the compact groups K_{∞} and K_{fin} each have total volume 1. This measure dg is independent of the choice of N since

$$
\operatorname{meas}(D_N)\operatorname{meas}(K_0(N)) = [\overline{\Gamma_0(1)} : \overline{\Gamma_0(N)}] \operatorname{meas}(D_1)\psi(N)^{-1} = \operatorname{meas}(D_1).
$$

Taking $N = 1$, a well-known computation of Siegel ([\[Si2\]](#page-69-9)) gives

$$
\operatorname{meas}(\mathbb{G}(\mathbf{Q})\backslash\mathbb{G}(\mathbf{A})) = \operatorname{meas}(\operatorname{Sp}_{2n}(\mathbf{Z})\backslash\mathcal{H}_n) = 2\prod_{j=1}^n [(j-1)!\pi^{-j}\zeta(2j)].
$$

Given $F \in S_{\mathbf{k}}(N)$, its adelic counterpart is the function $\varphi_F \in L_0^2$ defined by

$$
\varphi_F(g) = F(g_{\infty} \cdot iI_n)j(g_{\infty}, iI_n)^{-k}
$$
\n(7.8)

for $g = \gamma(g_{\infty} \times k) \in G(\mathbf{A}) = G(\mathbf{Q})(G(\mathbf{R})^+ \times K_0(N))$. The well-definedness of φ_F is a consequence of the fact that $G(\mathbf{Q}) \cap (G(\mathbf{R})^+ \times K_0(N)) = \Gamma_0(N)$. With measures normalized as above, the map $F \mapsto \varphi_F$ defines a linear isometry from $S_{\mathbf{k}}(N)$ onto the subspace $A_{\mathbf{k}}(N) \subseteq L_0^2$ defined in [\(4.2\)](#page-13-1). This may be proven just as in [\[Sa\]](#page-69-10). (The latter paper works with the principal congruence subgroup of level N , but the Siegel parabolic case is just the same.) The relationship between the adelic and classical Fourier coefficients is given by the following (see (6.2)).

PROPOSITION 7.1. For $F \in S_{\mathbf{k}}(N)$ and $\sigma \in S_n(\mathbf{Q}),$

$$
c_{\varphi_F}(\sigma) = \begin{cases} e^{-2\pi \operatorname{tr} \sigma} a_F(\sigma) & \text{if } \sigma \in \mathcal{R}_n^+ \\ 0 & \text{otherwise.} \end{cases}
$$
(7.9)

Proof. By strong approximation for the adeles,

$$
S_n(\mathbf{A}) = S_n(\mathbf{Q}) + S_n(\mathbf{R}) \times S_n(\mathbf{Z}).
$$

Decomposing $S = X_{\mathbf{Q}} + (X_{\infty} \times X_{\text{fin}}) \in S_n(\mathbf{A})$ accordingly, by the right $K_0(N)$ invariance and left $G(\mathbf{Q})$ -invariance of φ_F we have

$$
\varphi_F(\begin{pmatrix} I_n & S \\ & I_n \end{pmatrix}) = \varphi_F(\begin{pmatrix} I_n & X_\infty \\ & I_n \end{pmatrix}).
$$

Also, it follows that

$$
S_n(\mathbf{Q})\backslash S_n(\mathbf{A})=S_n(\mathbf{Z})\backslash (S_n(\mathbf{R})\times S_n(\mathbf{Z})).
$$

If D is any fundamental domain for $S_n(\mathbf{Z})\backslash S_n(\mathbf{R})$, then $D \times S_n(\widehat{\mathbf{Z}})$ is a fundamentaldomain for $S_n(\mathbf{Z}) \setminus (S_n(\mathbf{R}) \times S_n(\hat{\mathbf{Z}}))$ ([\[KL1,](#page-68-7) Theorem 7.40]). Therefore [\(6.2\)](#page-21-1) becomes

$$
c_{\varphi_F}(\sigma) = \int_{S_n(\mathbf{Z}) \backslash S_n(\mathbf{R})} \varphi_F(\begin{pmatrix} I_n & X_{\infty} \\ & I_n \end{pmatrix}) \theta_{\infty}(-\operatorname{tr}(\sigma X_{\infty})) dX_{\infty}
$$

$$
\times \int_{S_n(\widehat{\mathbf{Z}})} \theta_{\text{fin}}(-\operatorname{tr}(\sigma X_{\text{fin}})) dX_{\text{fin}}.
$$

Write $X_{fin} = (X_{ij})$ for $X_{ij} \in \mathbf{A}_{fin}$. If we write $\sigma = (b_{ij} \sigma_{ij})$, where b_{ij} is equal to 1 or 1/2 according to whether or not $i = j$, then $tr(\sigma X_{fin}) = \sum_{i,j} \sigma_{ij} X_{ij}$. Note that $\sigma_{ij} \in \mathbf{Z}$ for all i, j if and only if σ is half-integral, i.e., $\sigma \in \mathcal{R}_n$. If this condition does not hold, then the second integral vanishes.

Assuming $\sigma \in \mathcal{R}_n$, using [\(6.1\)](#page-20-3) and [\(7.8\)](#page-28-0) we have

$$
c_{\varphi_F}(\sigma) = \int_{S_n(\mathbf{Z}) \backslash S_n(\mathbf{R})} F(X + iI_n)e^{-2\pi i \operatorname{tr}(\sigma X)} dX
$$

= $e^{-2\pi \operatorname{tr} \sigma} \int_{S_n(\mathbf{Z}) \backslash S_n(\mathbf{R})} F(X + iI_n)e^{-2\pi i \operatorname{tr}(\sigma(X + iI_n))} dX$
= $e^{-2\pi \operatorname{tr} \sigma} a_F(\sigma).$

 \Box

Letting $\mathcal{B}_{k}(N)$ be an orthogonal basis for $S_{k}(N)$, in the special case $S = \emptyset$ and $r = 1$, our main formula [\(6.3\)](#page-24-1) becomes

$$
\lim_{N \to \infty} \frac{1}{\psi(N)} \sum_{F \in \mathcal{B}_{k}(N)} \frac{a_{\sigma_{1}}(F) \overline{a_{\sigma_{2}}(F)}}{\|F\|^{2}} \qquad (7.10)
$$
\n
$$
= \sum_{A \in \text{GL}_{n}(\mathbf{Z})/ \pm I_{n} \atop t_{A\sigma_{1}} A = \sigma_{2}} (\det A)^{k} d_{k} (4\pi)^{nk - n(n-1)/4} \frac{(\det \sigma_{1})^{k - (n+1)/2}}{\prod_{j=1}^{n} \Gamma(k - \frac{j-1}{2})}.
$$

For the formal degree, we will take the classical measure on $Sp(2n)$ corresponding to the measure in [\(7.7\)](#page-28-1). So as shown in [\[PSS\]](#page-69-7),

$$
d_{\mathbf{k}} = \frac{1}{2^{n} (4\pi)^{n(n+1)/2}} \prod_{1 \leq i \leq j \leq n} (2\mathbf{k} - (i+j)).
$$

(See the final remark after Proposition [A.7](#page-59-2) below.) We can simplify $\frac{d_k}{\Gamma_n(k)}$ using the following.

LEMMA 7.2. For any integers $n \geq 1$ and $k > 2n$,

$$
\frac{\prod_{1 \le i \le j \le n} (2k - (i + j))}{\prod_{\ell=1}^n \Gamma(\mathbf{k} - \frac{\ell - 1}{2})} = \frac{2^{n(n+1)/2}}{\prod_{j=1}^n \Gamma(\mathbf{k} - \frac{n+j}{2})}.
$$

Proof. When $n = 1$, by the functional equation for the Gamma function we have

$$
\frac{(2k-2)}{\Gamma(k)} = \frac{2(k-1)}{\Gamma(k)} = \frac{2}{\Gamma(k-1)},
$$

as needed. Given $n \geq 2$, suppose the formula holds for $n - 1$. Then

$$
\frac{\prod_{1 \le i \le j \le n} (2k - (i + j))}{\prod_{\ell=1}^n \Gamma(k - \frac{\ell - 1}{2})} = \frac{\prod_{1 \le i \le j \le n-1} (2k - (i + j))}{\prod_{\ell=1}^{n-1} \Gamma(k - \frac{\ell - 1}{2})} \frac{\prod_{1 \le i \le n} (2k - (i + n))}{\Gamma(k - \frac{n - 1}{2})}
$$
\n
$$
= \frac{2^{n(n-1)/2}}{\prod_{j=1}^{n-1} \Gamma(k - \frac{n-1+j}{2})} \frac{\prod_{1 \le i \le n} (2k - (i + n))}{\Gamma(k - \frac{n - 1}{2})}
$$
\n
$$
= 2^{n(n+1)/2} \prod_{i=1}^n \frac{(k - \frac{i + n}{2})}{\Gamma(k - \frac{i + n}{2} + 1)} = \frac{2^{n(n+1)/2}}{\prod_{i=1}^n \Gamma(k - \frac{n + i}{2})}.
$$

It now follows that

$$
\lim_{N \to \infty} \frac{1}{\psi(N)} \sum_{F \in \mathcal{B}_k(N)} \frac{a_{\sigma_1}(F) \overline{a_{\sigma_2}(F)}}{\|F\|^2} \n= \frac{(\det \sigma_1)^{k - (n+1)/2}}{\pi^{n(n-1)/4} (4\pi)^{n(n+1)/2 - nk} \prod_{j=1}^n \Gamma(k - \frac{n+j}{2})} \delta_k(\sigma_1, \sigma_2),
$$
\n(7.11)

for $\delta_{\mathbf{k}}(\sigma_1, \sigma_2)$ given in [\(1.3\)](#page-4-0). This proves Theorem [1.1.](#page-3-0)

7.4 NONARCHIMEDEAN INTEGRALS: $p \in \mathbf{S}$

This case is more difficult. Our goal is to compute (or bound) the local integral

$$
I_{A,p} = I_{A,p}(f_p) = \int_{S_n(\mathbf{Q}_p)} f_p(\begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix} \begin{pmatrix} A & O \\ O & r^t A^{-1} \end{pmatrix}) \theta_p(\text{tr }\sigma_1 S) dS.
$$

To simplify the computation, we may essentially reduce to the case where A is diagonal, as follows. By the elementary divisors theorem, there exist $U, V \in GL_n(\mathbf{Z})$ and a diagonal matrix

$$
D = \text{diag}(d_1, \ldots, d_n)
$$

with positive integer entries satisfying $d_1|d_2|\cdots|d_n$, such that

$$
A = UDV.\t(7.12)
$$

PROPOSITION 7.3. For A as above,

$$
I_{A,p} = \int_{S_n(\mathbf{Q}_p)} f_p(\begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix} \begin{pmatrix} D & O \\ O & rD^{-1} \end{pmatrix}) \theta_p(\text{tr }\sigma'_1 S) dS,
$$

where

$$
\sigma_1' = \sigma_U = {}^tU\sigma_1U.
$$

Proof. By definition,

$$
I_{A,p} = \int_{S_n(\mathbf{Q}_p)} f_p(\begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix} \begin{pmatrix} A & O \\ O & r^t A^{-1} \end{pmatrix}) \theta_p(\text{tr }\sigma_1 S) dS
$$

=
$$
\int_{S_n(\mathbf{Q}_p)} f_p(\begin{pmatrix} UDV & rS^tU^{-1}D^{-1}tV^{-1} \\ O & r^tU^{-1}D^{-1}tV^{-1} \end{pmatrix}) \theta_p(\text{tr }\sigma_1 S) dS.
$$

Because f_p is bi- K_p -invariant, we are free to multiply its argument on the left by $\int U^{-1}$ tU $\Big) \in K_p$ and on the right by $\left(\begin{matrix} V^{-1} \end{matrix}\right)$ \sqrt{t} \setminus $\in K_p$. This gives

$$
I_{A,p} = \int_{S_n(\mathbf{Q}_p)} f_p(\begin{pmatrix} D & r(U^{-1}S^t U^{-1})D^{-1} \\ O & rD^{-1} \end{pmatrix}) \theta_p(\text{tr }\sigma_1 S) dS.
$$

Let $S' = U^{-1}S^tU^{-1}$. Then $dS' = dS$ since $S \mapsto S'$ is an isomorphism mapping $S(\widehat{\mathbf{Z}})$ to $S(\widehat{\mathbf{Z}})$. Hence the above is

$$
= \int_{S_n(\mathbf{Q}_p)} f_p(\begin{pmatrix} D & rS'D^{-1} \\ O & rD^{-1} \end{pmatrix}) \theta_p(\text{tr }\sigma_1 US' {}^tU) dS'.
$$

Now using $\operatorname{tr} \sigma_1 US^tU = \operatorname{tr}^tU\sigma_1 US$, we find

$$
I_{A,p} = \int_{S_n(\mathbf{Q}_p)} f_p(\begin{pmatrix} D & rSD^{-1} \\ O & rD^{-1} \end{pmatrix}) \theta_p(\text{tr }\sigma'_1 S) dS.
$$

For the purpose of computing the above local integral, by the K_p -invariance of f_p , we may assume that each diagonal entry of D is a power of p. Thus, we take

$$
D = \text{diag}(p^{\alpha_1}, \dots, p^{\alpha_n}),
$$

for

$$
0\leq \alpha_1\leq \cdots \leq \alpha_n.
$$

For $x \in \mathbf{Q}_p$, we define $\text{ord}_p(x) = n$ if $x = p^n z$ with z a unit. So $r_p = \text{ord}_p(r)$, for example.

PROPOSITION 7.4. With the above notation, suppose $I_{A,p} \neq 0$. Then each $\alpha_j \leq r_p$. Under the additional assumption that $\text{ord}_p(\det \sigma_1) = \text{ord}_p(\det \sigma_2)$, we further have

$$
\alpha_1 + \cdots + \alpha_n = \frac{nr_p}{2}.
$$

Proof. As in Theorem [6.3,](#page-24-1) we are assuming that $r^t A^{-1} \in M_n(\mathbf{Z})$. It follows that likewise $rD^{-1} \in M_n(\mathbf{Z})$, and hence $\alpha_j \leq r_p$ for each j. Under the additional assumption, taking determinants in the relation ${}^t\!A\sigma_1A = r\sigma_2$ gives $p^{2(\alpha_1+\cdots+\alpha_n)}=p^{nr_p}$, and the last assertion follows. \Box

In principle, one can now compute the integral by applying Proposition [3.2](#page-12-1) and considering various cases to obtain certain exponential sums. We will discuss this process in more detail for the special case of GSp(4) in Section [9.](#page-40-0) It should be evident from this special case that the general case is very complicated.

We conclude the present section by giving a trivial bound for $I_{A,p}$.

PROPOSITION 7.5. With notation as above,

$$
|I_{A,p}| \le \prod_{j=1}^{n} p^{j(r_p - \alpha_j)} = p^{\frac{n(n+1)}{2}r_p - (\alpha_1 + 2\alpha_2 + \dots + n\alpha_n)}.
$$
 (7.13)

Proof. By Proposition [5.1,](#page-15-1)

$$
f_p(\begin{pmatrix} D & rSD^{-1} \\ O & rD^{-1} \end{pmatrix}) \neq 0 \iff \begin{pmatrix} D & rSD^{-1} \\ O & rD^{-1} \end{pmatrix} \in C_p.
$$

Because $C_p \subseteq M_{2n}(\mathbf{Z}_p)$, we see that

$$
|I_{A,p}| \leq \text{meas}\{S \in S_n(\mathbf{Q}_p) | p^{r_p} S D^{-1} \in M_n(\mathbf{Z}_p)\}.
$$

Writing $S = (s_{ij}),$

$$
|I_{A,p}| \leq \text{meas}\{S \in S_n(\mathbf{Q}_p) | s_{ij}p^{r_p - \alpha_j} \in \mathbf{Z}_p \text{ for all } i, j\}
$$

=
$$
\text{meas}\{S \in S_n(\mathbf{Q}_p) | s_{ij} \in p^{-(r_p - \alpha_j)}\mathbf{Z}_p \text{ for all } 1 \leq i \leq j \leq n\}
$$

since $s_{ij} = s_{ji}$. Hence,

$$
|I_{A,p}| \le \prod_{j=1}^n \prod_{i \le j} \text{meas}(p^{-(r_p - \alpha_j)} \mathbf{Z}_p) = \prod_{j=1}^n \prod_{i \le j} p^{r_p - \alpha_j} = \prod_{j=1}^n p^{j(r_p - \alpha_j)}.
$$

710 Andrew Knightly and Charles Li

8 Weighted equidistribution of Satake parameters

Let $\hat{\mathbb{G}} = \text{Spin}(2n + 1, \mathbb{C})$ be the complex dual group of $\mathbb{G} = \text{PGSp}_{2n}$.^{[†](#page-33-3)} Since we are assuming trivial central character, the Satake parameters t_{π_n} belong to the maximal torus $\hat{\mathbb{T}}$ of $\hat{\mathbb{G}}$. For our fixed finite set **S** of primes, let

$$
\mathfrak{X}_{\mathbf{S}} = (\widehat{\mathbb{T}}/W)^{|\mathbf{S}|}.
$$

Let $\Pi_{\mathbf{k}}(N)$ be the set of cuspidal representations of weight k and level N defined in §[4.](#page-12-0) Each $\pi \in \Pi_{\mathbf{k}}(N)$ determines a point

$$
t_{\pi} = (t_{\pi_p})_{p \in \mathbf{S}} \in \mathfrak{X}_{\mathbf{S}}.\tag{8.1}
$$

 \Box

By Shin's theorem, the points t_{π} become equidistributed relative to the Plancherel product measure on $\mathfrak{X}_\mathbf{S}$ as $N \to \infty$. Here we investigate their distribution with certain prescribed harmonic weights.

PROPOSITION 8.1. There exists a compact subset $\Omega \subseteq \mathfrak{X}_{\mathbf{S}}$ such that $t_{\pi} \in \Omega$ for all $\pi \in \Pi_{\mathbf{k}}(N)$.

Proof. [\[BW,](#page-67-10) Theorem XI.3.3].

8.1 Preliminary result

For $\sigma_1, \sigma_2 \in \mathcal{R}_n^+$, define the weight

$$
w_{\pi}(\sigma_1, \sigma_2) = \sum_{\varphi \in E_{\mathsf{k}}(\pi, N)} \frac{c_{\sigma_1}(\varphi) c_{\sigma_2}(\varphi)}{\|\varphi\|^2}.
$$

We will show in this section that the Satake parameters t_{π} , weighted by $w_{\pi}(\sigma,\sigma)$, have a uniform distribution relative to a certain Radon measure in the limit as $N \to \infty$.

The following is essentially a restatement of Theorem [1.1](#page-3-0) (see §[7.3\)](#page-27-1).

LEMMA 8.2. Let
$$
c_{nk\sigma_1} = \frac{(\det \sigma_1)^{k-(n+1)/2}}{\pi^{n(n-1)/4} (4\pi)^{n(n+1)/2 - nk} \prod_{j=1}^n \Gamma(k - \frac{n+j}{2})}
$$
. Then

$$
\lim_{N \to \infty} \frac{1}{\psi(N)} \sum_{\pi \in \Pi_k(N)} w_{\pi}(\sigma_1, \sigma_2) = \delta_k(\sigma_1, \sigma_2) c_{nk\sigma_1},
$$

where $\delta_{\mathbf{k}}(\sigma_1, \sigma_2)$ is defined in [\(1.3\)](#page-4-0). In particular, if $\sigma_1 = \sigma_2 = \sigma \in \mathcal{R}_n^+$ satisfies $\delta_{\bf k}(\sigma,\sigma) > 0$ (e.g. if **k** is even), then setting $w_{\pi} = w_{\pi}(\sigma,\sigma)$ we have

$$
0 < \lim_{N \to \infty} \frac{1}{\psi(N)} \sum_{\pi \in \Pi_k(N)} w_{\pi} < \infty. \tag{8.2}
$$

†For dual groups, we always take the ground field to be C unless specified otherwise.

Remark: If k is even, then $\delta_{\mathbf{k}}(\sigma,\sigma) > 0$ since $A = I_n$ satisfies ${}^t A \sigma A = \sigma$. Hence in this case we always have nonvanishing in (8.2) . This nonvanishing is crucial in what follows.

For the compact space Ω in Proposition [8.1,](#page-33-5) let

$$
V_{\mathcal{S}} \subseteq C(\Omega) \tag{8.3}
$$

denote the subspace consisting of all restrictions $F|_{\Omega}$ of functions $F = \prod_{p \in S} S f_p$ in the image $\prod_{p \in S} C[X^*(\hat{T})]^W$ of the S-product Satake transform. (See [\(2.3\)](#page-6-5) and [\(2.4\)](#page-6-1).) For fixed $k > 2n$ and $\sigma \in \mathcal{R}_n^+$ for which $\delta_k(\sigma, \sigma) > 0$, let $w_{\pi} = w_{\pi}(\sigma, \sigma)$ as above. Then we may define a linear functional $\mathcal{L} = \mathcal{L}_{\sigma, k, S}$ on V_S by

$$
\mathcal{L}(F) = \lim_{\substack{N \to \infty \\ (N, S) = 1}} \frac{\sum_{\pi \in \Pi_{k}(N)} w_{\pi} F(t_{\pi})}{\sum_{\pi \in \Pi_{k}(N)} w_{\pi}}.
$$
\n(8.4)

By Lemma [8.2](#page-33-6) and Theorem [6.3,](#page-24-1) the limit exists and is finite. Endowing $C(\Omega)$ with the L^{∞} norm, it contains $V_{\mathcal{S}}$ as a dense subalgebra by the Stone-Weierstrass Theorem [\[Ru,](#page-69-11) p. 122]. (The latter algebra evidently separates points, and it is closed under complex conjugation by Proposition [2.1.](#page-7-4)) Because V_S is dense in $C(\Omega)$, the right-hand side of [\(8.4\)](#page-34-1) exists for $F \in C(\Omega)$ (for details, see e.g. [\[KL1\]](#page-68-7), pages 358-359). Moreover it is clear from [\(8.4\)](#page-34-1) that $|\mathcal{L}(F)| \leq ||F||_{\infty}$ for all $F \in C(\Omega)$, so $\mathcal L$ is bounded. By the Riesz representation theorem, there exists a unique Radon measure $\mu = \mu_{\sigma,k,\mathbf{S}}$ on Ω such that

$$
\mathcal{L}(F) = \int_{\Omega} F d\mu \tag{8.5}
$$

for all $F \in C(\Omega)$. It is clear from [\(8.4\)](#page-34-1) that $\mathcal{L}(1) = 1$, so μ is a probability measure. This proves the following.

THEOREM 8.3. Let $k > 2n$, and let σ be a symmetric positive-definite halfintegral matrix for which $\delta_{\mathbf{k}}(\sigma,\sigma) > 0$. (This is automatic, for example, if \mathbf{k} is even.) Then the Satake parameters $(t_{\pi})_{\pi \in \Pi_{k}(N)}$ of (8.1) , when weighted by $w_{\pi}(\sigma,\sigma)$, become equidistributed in the space Ω of Proposition [8.1](#page-33-5) with respect to the above probability measure $\mu_{\sigma,k,S}$ in the limit as $N \to \infty$ along integers coprime to S.

Of course, one would like to know more about the measure μ , for example whether it is supported on the tempered spectrum. (Recall that an unramified representation π_p of $G(\mathbf{Q}_p)$ is tempered if and only if its Satake parameter t_{π_p} lies in a compact subgroup of \overline{T} .) We will pursue this question by relating μ to the Sato-Tate measure, which is supported in a compact subtorus of $\hat{\mathbb{T}}$. See Theorem [8.4](#page-36-1) below.

712 Andrew Knightly and Charles Li

8.2 Relating two measures

Generally, suppose η is a Radon measure on Ω , and $\{R_{\lambda}\}_{\lambda \in \Lambda}$ is a set of continuous functions forming an orthonormal basis for $L^2(\Omega, \eta)$ which also spans an L[∞]-dense subspace of $C(\Omega)$. Then the measure μ in [\(8.5\)](#page-34-2) can be expressed as

$$
d\mu(t) = \sum_{\lambda} \mathcal{L}(R_{\lambda}) \overline{R_{\lambda}}(t) d\eta(t), \qquad (8.6)
$$

provided the sum is uniformly absolutely convergent on Ω . Indeed, for all $\alpha \in \Lambda$,

$$
\int_{\Omega} R_{\alpha}(t) \sum_{\lambda} \mathcal{L}(R_{\lambda}) \overline{R_{\lambda}}(t) d\eta = \sum_{\lambda} \mathcal{L}(R_{\lambda}) \langle R_{\alpha}, R_{\lambda} \rangle_{\eta} = \mathcal{L}(R_{\alpha}) = \int_{\Omega} R_{\alpha} d\mu,
$$

and by linearity and density of the span of the R_{α} , the above holds as well for all functions in $C(\Omega)$.

8.3 THE SATO-TATE MEASURE

Fix a maximal compact subgroup $\widehat{H} \subseteq \widehat{\mathbb{G}}$ with maximal torus

$$
\widehat{\mathbb{T}}_c = \widehat{H} \cap \widehat{\mathbb{T}} = \widehat{\mathbb{T}}(\mathbf{C})^1.
$$

This is the maximal compact subtorus of $\widehat{\mathbb{T}}$. Let dh denote the Haar measure on H of total volume 1. Because every conjugacy class in H contains exactly one Weyl orbit of \hat{T}_c , the measure dh induces a quotient measure μ_{ST} on the space $\widehat{\mathbb{T}}_c/W$. We extend μ_{ST} to $\widehat{\mathbb{T}}/W$ by taking it to be zero on the complement of $\widehat{\mathbb{T}}_c/W$. This is the Sato-Tate measure. In more detail, for $f \in C(\widehat{\mathbb{T}}_c/W)$, we may identify f with a class function on \hat{H} , and

$$
\int_{\widehat{\mathbb{T}}_c/W} f(t) d\mu_{ST}(t) = \int_{\widehat{H}} f(h) dh.
$$
\n(8.7)

By the Weyl integration formula, the measure is given explicitly by

$$
d\mu_{ST}(t) = \left| \det(\mathrm{Ad}(t^{-1} - I)|_{\mathrm{Lie}(\widehat{H})/\mathrm{Lie}(\widehat{T}_c)}) \right| dt, \tag{8.8}
$$

where dt is the Haar measure giving \hat{T}_c volume 1. An alternative expression for it is given in [\(8.13\)](#page-37-0) below.

Fix a set Φ^+ of positive roots in the root system attached to \hat{H} and \hat{T}_c . We shall identify $X^*(\hat{\mathbb{T}}_c)$ and $X^*(\hat{\mathbb{T}}_c)$. By the theorem of the highest weight, the irreducible representations π_{λ} of \widehat{H} are in one-to-one correspondence with the elements $\lambda \in C^+$, where C^+ is the positive Weyl chamber of $X^*(\mathbb{T}_c) = X^*(\mathbb{T}) \cong$ $X_*(\mathbb{T})$ given in [\(3.19\)](#page-11-2). Let

$$
F_{\lambda} = \operatorname{tr} \pi_{\lambda} \tag{8.9}
$$
denote the trace of π_{λ} . It is a class function on \widehat{H} , so we may view it as a function on $\widehat{\mathbb{T}}_c/W$. By the Peter-Weyl theorem, the set $\{F_\lambda | \lambda \in C^+\}$ is an orthonormal basis for the space of L^2 class functions on \hat{H} (relative to the measure dh). In particular, by (8.7) we have

$$
\int_{\widehat{\mathbb{T}}(\mathbf{C})/W} F_{\lambda}(t) \overline{F_{\mu}(t)} d\mu_{ST} = \delta_{\lambda,\mu}
$$
\n(8.10)

(for the Kronecker δ). Here, the domain of F_{λ} is extended from $\widehat{\mathbb{T}}_c$ to $\widehat{\mathbb{T}}(\mathbf{C})$ by viewing F_{λ} as a sum

$$
F_{\lambda} = \sum_{\mu \in X^*(\widehat{\mathbb{T}})} m_{\lambda}(\mu)[\mu] \in \mathbf{C}[X^*(\widehat{\mathbb{T}})]^W.
$$

The orthogonality [\(8.10\)](#page-36-0) can also be proved using [\(8.12\)](#page-37-0) and [\(8.13\)](#page-37-1) below.

We shall need the fact that the set $\{F_{\lambda} | \lambda \in C^{+}\}\$ spans $\mathbb{C}[X^*(\mathbb{T})]^W$ (see [\[FH,](#page-67-0) Theorem 23.24], using the fact that $\Lambda = X^*(\mathbb{T})$ since $\mathbb{G} = \text{Spin}(2n + 1)$ is simply connected). By (2.4) , this space coincides with the image of the local Satake transform.

Given a tuple $\underline{\lambda} = (\lambda_p) \in \prod_{p \in \mathbf{S}} C^+$, we let

$$
F_{\underline{\lambda}} = \prod_{p \in \mathbf{S}} F_{\lambda_p} \in \prod_{p \in \mathbf{S}} \mathbf{C}[X^*(\widehat{\mathbb{T}})]^W.
$$

Viewing the $F_{\underline{\lambda}}$ as functions on Ω (by restriction), they span the space $V_{\mathcal{S}}$ of [\(8.3\)](#page-34-0). This follows from the above discussion.

8.4 RELATION BETWEEN μ and μ_{ST}

We continue to assume that $\delta_{\mathbf{k}}(\sigma,\sigma) > 0$, so that for $F_{\underline{\lambda}}$ as above, we may consider $\mathcal{L}(F_{\lambda})$ as in [\(8.4\)](#page-34-1).

THEOREM 8.4. Let $\rho \in X^*(\mathbb{T}) = X_*(\mathbb{T})$ be half the sum of the positive roots, as in [\(3.10\)](#page-9-0). Suppose that there exits $\varepsilon > 0$ such that for all tuples λ as above,

$$
\mathcal{L}(\prod_{p \in \mathbf{S}} \mathcal{S}(c_{\lambda_p})) \ll_{\varepsilon} \prod_{p \in \mathbf{S}} p^{(1-\varepsilon)\langle \rho, \lambda_p \rangle},\tag{8.11}
$$

where c_{λ_p} is the characteristic function of $K_p\lambda_p(p)K_p$. Then the measure μ defined in [\(8.5\)](#page-34-2) is given by

$$
d\mu(t) = \sum_{\underline{\lambda}} \mathcal{L}(F_{\underline{\lambda}}) \overline{F_{\underline{\lambda}}(t)} d\mu \mathbf{s}(t),
$$

where $\mu_{\mathbf{S}} = \prod_{p \in \mathbf{S}} \mu_{ST}$ is the product measure on $\mathfrak{X}_{\mathbf{S}}$, and the above sum converges absolutely and uniformly on Ω .

Remarks: (1) Hypothesis [\(8.11\)](#page-36-1) would follow from (a) adequate bounds on the number of matrices A satisfying the conditions of Theorem [6.3,](#page-24-0) and (b) adequate bounds for the local geometric integrals $I_{A,p}(c_{\lambda_p})$ for $p \in S$. (See [\(9.3\)](#page-40-0).) In Section [9,](#page-40-1) we will carry this out and prove Hypothesis [\(8.11\)](#page-36-1) in the special case $n = 2$, as an application of Theorem [6.3.](#page-24-0)

(2) It is not clear to us whether there is a closed form expression for the measure.

Proof. In Section [8.3,](#page-35-1) we saw that the set $\{F_{\lambda} | \lambda \in C^{+}\}\$ is an orthonormal basis for $L^2(\Omega, \mu_{ST})$. Furthermore, it is dense in $V_{\mathcal{S}}$, which in turn is dense in $C(\Omega)$ as discussed before Theorem [8.3.](#page-34-3) Hence, by the discussion in §[8.2,](#page-35-2) it suffices to prove that the given series is uniformly convergent under Hypothesis [\(8.11\)](#page-36-1). To ease the notation, we will first assume that $S = \{p\}$ consists of just one prime. For any weight $\lambda \in X^*(\hat{\mathbb{T}})$, define

$$
A_{\lambda} = \sum_{w \in W} (\operatorname{sgn} w) w(\lambda) \in \mathbf{C}[X^*(\widehat{\mathbb{T}})].
$$

For $t \in \widehat{\mathbb{T}}_c$,

$$
|A_{\lambda}(t)| \leq \left| \sum_{w \in W} (\operatorname{sgn} w) w(\lambda)(t) \right| \leq \sum_{w \in W} |w(\lambda)(t)| = |W|.
$$

By the Weyl character formula([\[FH,](#page-67-0) Theorem 24.2]),

$$
F_{\lambda} = \frac{A_{\lambda+\rho}}{A_{\rho}}.\tag{8.12}
$$

It is well-known that

$$
d\mu_{ST}(t) = |A_{\rho}(t)|^2 dt. \tag{8.13}
$$

(For example, compare (25.6) of [\[Bu\]](#page-67-1) ((22.7) in the 2nd edition) with Lemma 24.3 of [\[FH\]](#page-67-0)).

Therefore, for $t \in \hat{T}_c$, we need to prove the convergence of

$$
\sum_{\lambda \in \mathcal{C}^+} |\mathcal{L}(F_{\lambda}) \overline{F_{\lambda}(t)}| |A_{\rho}(t)|^2 = \sum_{\lambda \in \mathcal{C}^+} |\mathcal{L}(F_{\lambda}) \overline{A_{\lambda+\rho}(t)} A_{\rho}(t)| \le |W|^2 \sum_{\lambda \in \mathcal{C}^+} |\mathcal{L}(F_{\lambda})|.
$$

Next, we need to relate F_{λ} to the functions $\mathcal{S}(c_{\mu})$ in order to make use of Hypothesis [\(8.11\)](#page-36-1). This is achieved by the following formula of Kato and Lusztig,which holds in any split reductive p -adic group ([\[HKP,](#page-68-0) Theorem 7.8.1]; see also [\[Gr,](#page-67-2) (3.12) and Proposition 4.4]):

$$
F_{\lambda} = p^{-\langle \lambda, \rho \rangle} \sum_{\mu \leq \lambda} P_{\mu, \lambda}(p) \mathcal{S}(c_{\mu}).
$$

Here, μ belongs to C^+ , and $P_{\mu,\lambda}$ is the Kazhdan-Lusztig polynomial

$$
P_{\mu,\lambda}(p) = p^{\langle \lambda - \mu, \rho \rangle} \sum_{w \in W} \text{sgn}(w) \widehat{P}(w(\lambda + \rho^{\vee}) - (\mu + \rho^{\vee})),
$$

where

$$
\widehat{P}(\mu) = \sum_{\mu = \sum n(\alpha^{\vee})\alpha^{\vee}} p^{-\sum n(\alpha^{\vee})} \ge 0
$$

encodes the number of expressions of μ as a linear combination of positive co-roots with coefficients $n(\alpha^{\vee}) \ge 0$. We note that $P_{\lambda,\lambda}(p) = 1$, [\[Gr,](#page-67-2) (4.5)]. Therefore, the quantity we need to bound is

$$
\sum_{\lambda \in \mathcal{C}^+} |\mathcal{L}(F_{\lambda})| = \sum_{\lambda \in \mathcal{C}^+} \left| p^{-\langle \lambda, \rho \rangle} \sum_{\mu \le \lambda} P_{\mu, \lambda}(p) \mathcal{L}(\mathcal{S}(c_{\mu})) \right|
$$

$$
\le \sum_{\mu \in \mathcal{C}^+} p^{-\langle \mu, \rho \rangle} |\mathcal{L}(\mathcal{S}(c_{\mu}))| \sum_{\lambda \ge \mu} \sum_{w \in W} \widehat{P}(w(\lambda + \rho^{\vee}) - (\mu + \rho^{\vee})). \tag{8.14}
$$

We claim that for $w \in W$,

$$
\sum_{\lambda \ge \mu} \widehat{P}(w(\lambda + \rho^{\vee}) - (\mu + \rho^{\vee})) \le 2^{d^+},\tag{8.15}
$$

where d^+ is the number of positive co-roots. Indeed, the left-hand side of (8.15) is

$$
\sum_{\lambda \ge \mu} \sum_{\sum n(\alpha^{\vee}) \alpha^{\vee}} p^{-\sum n(\alpha^{\vee})}, \tag{8.16}
$$

where the inner sum is extended over all expressions of the form

$$
w(\lambda + \rho^{\vee}) - (\mu + \rho^{\vee}) = \sum_{\alpha^{\vee}} n(\alpha^{\vee})\alpha^{\vee}
$$

with $n(\alpha^{\vee}) \geq 0$ and α^{\vee} positive co-roots. The above expression is equivalent to

$$
\lambda = w^{-1} \left(\sum_{\alpha^{\vee}} n(\alpha^{\vee}) \alpha^{\vee} + (\mu + \rho^{\vee}) \right) - \rho^{\vee}.
$$
 (8.17)

Thus we may exchange the order of summation in [\(8.16\)](#page-38-1), so the left-hand side of [\(8.15\)](#page-38-0) is equal to

$$
\sum_{\sum n(\alpha^{\vee})\alpha^{\vee}}^* p^{-\sum n(\alpha^{\vee})},
$$

where the ∗ indicates that we consider only those expressions for which the right-hand side of [\(8.17\)](#page-38-2) is $\geq \mu$. The above is of course bounded by the sum over all nonnegative linear combinations of positive co-roots

$$
\sum_{\sum n(\alpha^{\vee})\alpha^{\vee}} p^{-\sum n(\alpha^{\vee})} = \prod_{\alpha^{\vee}} \sum_{n(\alpha^{\vee})=0}^{\infty} p^{-n(\alpha^{\vee})} \le \prod_{\alpha^{\vee}} 2,
$$

proving the claim [\(8.15\)](#page-38-0).

Combining [\(8.14\)](#page-38-3) and [\(8.15\)](#page-38-0), it follows that

$$
\sum_{\lambda \in \mathcal{C}^+} |\mathcal{L}(F_\lambda)| \ll |W| \sum_{\mu \in \mathcal{C}^+} p^{-\langle \mu, \rho \rangle} |\mathcal{L}(\mathcal{S}(c_\mu))|.
$$

Using the given bound [\(8.11\)](#page-36-1), the above is

$$
\ll \sum_{\mu \in \mathcal{C}^+} p^{-\varepsilon \langle \mu, \rho \rangle}.
$$

There exists a finite set $\{\mu_1,\ldots,\mu_\ell\} \subseteq C^+$ such that $C^+ \subseteq {\sum_{i=1}^\ell a_i \mu_i | 0 \leq \ell\}}$ $a_i \in \mathbf{Z}$. Writing $\mu = \sum_{i=1}^{\ell} a_i \mu_i$, the above is

$$
\leq \prod_{i=1}^{\ell} \left(\sum_{a_i=0}^{\infty} p^{-\varepsilon \langle \mu_i, \rho \rangle a_i} \right) < \infty.
$$

This completes the proof when $S = \{p\}.$

The general case is proven in the same way, using

$$
\sum_{\Delta} |\mathcal{L}(F_{\Delta})| = \sum_{\Delta} \left| \sum_{\mu \leq \Delta} \left(\prod_{p \in \mathbf{S}} p^{-\langle \lambda_p, \rho \rangle} P_{\mu_p, \lambda_p}(p) \right) \mathcal{L}(\prod_{p \in \mathbf{S}} \mathcal{S}(c_{\mu_p})) \right|
$$

$$
\leq \sum_{\mu} |\mathcal{L}(\prod_{p \in \mathbf{S}} \mathcal{S}(c_{\mu_p}))| \prod_{p \in \mathbf{S}} \sum_{\lambda_p \geq \mu_p} p^{-\langle \lambda_p, \rho \rangle} |P_{\mu_p, \lambda_p}(p)|
$$

$$
\ll \sum_{\mu} |\mathcal{L}(\prod_{p \in \mathbf{S}} \mathcal{S}(c_{\mu_p}))| \prod_{p \in \mathbf{S}} p^{-\langle \mu_p, \rho \rangle}.
$$

Using Hypothesis [8.11,](#page-36-1) one shows as before that this is finite.

COROLLARY 8.5. Write $\mu = \mu_p$ for the measure on \hat{T}/W defined in [\(8.5\)](#page-34-2) when $S = \{p\}$. Then under Hypothesis [8.11,](#page-36-1)

$$
\lim_{p \to \infty} d\mu_p(t) = d\mu_{ST}(t).
$$

Proof. Let $0 \in C^+$ denote the element corresponding to the zero vector in \mathbf{Z}^{n+1} . As in the proof of the previous proposition,

$$
\sum_{\lambda \in \mathcal{C}^+ - \{\mathbf{0}\}} |\mathcal{L}(F_\lambda) \overline{F_\lambda(t)}| |A_\rho(t)|^2 \ll \sum_{\mu \in \mathcal{C}^+ - \{\mathbf{0}\}} p^{-\varepsilon \langle \mu, \rho \rangle}.
$$

Noting that $\langle \mu, \rho \rangle > 0$ when $\mu \neq \mathbf{0}$, the right-hand side tends to 0 as p goes to ∞. Thus

$$
\lim_{p \to \infty} d\mu_p(t) = \lim_{p \to \infty} \sum_{\lambda \in \mathcal{C}^+} \mathcal{L}(F_\lambda) \overline{F_\lambda(t)} d\mu_{ST}(t) = \mathcal{L}(F_0) F_0(t) d\mu_{ST}(t) = d\mu_{ST}(t).
$$

The last step follows by [\(8.4\)](#page-34-1) and the fact that $F_0 = 1$ (cf. [\(8.12\)](#page-37-0)).

 \Box

 \Box

9 LOCAL COMPUTATION WHEN $n = 2$

Here we refine the discussion from Section [7.4](#page-30-0) for $p \in S$, with the simplifying assumptions that $n = 2$, and

$$
p \nmid 4 \det \sigma_1. \tag{9.1}
$$

(Recall that det $\sigma_1 \in \frac{1}{4}\mathbb{Z}$.) The main goal of this section is to prove the following local bound.

PROPOSITION 9.1. Under the above hypotheses, there exists a constant $\varepsilon > 0$ such that

$$
|I_{A,p}(c_{\lambda})| \ll p^{(1-\varepsilon)\langle \lambda,\rho\rangle - \varepsilon r_{p}}
$$
\n(9.2)

for all $\lambda \in C^+$, where the implied constant depends only on p and ε .

Before proving the proposition, let us observe how it implies the global Hypothesis [\(8.11\)](#page-36-1).

COROLLARY 9.2. Suppose $n = 2$, $\sigma_1 = \sigma_2 = \sigma$, $\delta_k(\sigma, \sigma) > 0$, and $p \nmid 4 \det \sigma$ for all $p \in \mathbf{S}$. Then Hypothesis [\(8.11\)](#page-36-1) holds.

Proof. In Section [7.2,](#page-27-0) we saw that $I_{A,p} = 1$ for primes $p \notin S$. From the definition (8.4) of \mathcal{L} , Theorem [6.3,](#page-24-0) and Lemma [8.2,](#page-33-0) it follows that

$$
\mathcal{L}(\prod_{p \in \mathbf{S}} \mathcal{S}(c_{\lambda_p})) = \frac{\sum_{A} c_{n \mathbf{k} \sigma} \prod_{p \in \mathbf{S}} I_{A,p}(c_{\lambda_p})}{\delta_{\mathbf{k}}(\sigma, \sigma) c_{n \mathbf{k} \sigma}} = \frac{1}{\delta_{\mathbf{k}}(\sigma, \sigma)} \sum_{A} \prod_{p \in \mathbf{S}} I_{A,p}(c_{\lambda_p}), \quad (9.3)
$$

where A runs through the rank-2 matrices in $M_2(\mathbf{Z})/\{\pm 1\}$ satisfying $r^tA^{-1} \in$ $M_2(\mathbf{Z})$ and ${}^t\!A\sigma A = r\sigma$. In particular, writing

$$
\sigma = \begin{pmatrix} a & b/2 \\ b/2 & c \end{pmatrix}, \qquad A = \begin{pmatrix} x & y \\ z & w \end{pmatrix} \in M_2(\mathbf{Z}),
$$

we have $ax^2 + bxz + cz^2 = ra$. Hence

$$
4ra2 = (2ax + bz)2 - (b2 - 4ac)z2 = (2ax + bz + \sqrt{D}z)(2ax + bz - \sqrt{D}z),
$$

where $D = b^2 - 4ac < 0$. Thus, in the ring of integers $\mathcal{O} \subseteq \mathbf{Q}[\sqrt{D}]$, the ideal $(2ax + bz + \sqrt{D}z)$ is a factor of the ideal $(4ra^2)$. The number of ideal factors of $(4ra^2)$ is $\ll r^{\epsilon/2}$. In view of the fact that $|\mathcal{O}^*| < \infty$, the number of possible choices for x, z is $\ll r^{\varepsilon/2}$. Similarly, the number of choices for y, w is $\ll r^{\varepsilon/2}$. So the number of terms in the sum is $\ll r^{\varepsilon} = \prod_{p \in \mathbf{S}} p^{r_p \varepsilon}$. It follows from [\(9.2\)](#page-40-2) that the above is

$$
\ll \prod_{p\in \mathbf{S}} p^{(1-\varepsilon)\langle \lambda_p,\rho\rangle},
$$

as required.

718 Andrew Knightly and Charles Li

The proof of Proposition [9.1](#page-40-3) is given in Section [9.3.](#page-49-0) In the intervening sections, we describe how to compute the local integral $I_{A,p}$ explicitly, with the goal of producing the upper bound [\(9.2\)](#page-40-2). In many situations, the trivial bound [\(7.13\)](#page-32-0) is adequate, so an explicit computation is not necessary. In the remaining cases (which, in the notation below, occur when $\beta - 1 \leq t$), we give a complete treatment of the local integral.

9.1 Preliminaries

Without loss of generality, we consider the case where $f_p = c_\lambda$ is the characteristic function of the double coset $Z_pK_p\lambda(p)K_p$, where

$$
\lambda(p) = \text{diag}(1, p^t, p^\tau, p^{\tau - t}) \tag{9.4}
$$

for $0 \le t \le \tau/2$ as in [\(3.21\)](#page-11-0). Thus we write τ in place of r_p (for the purpose of eliminating a subscript).

By Proposition [7.4,](#page-32-1) we need only consider matrices D of the form

$$
D = \text{diag}(p^{\alpha}, p^{\beta}), \quad 0 \le \alpha \le \beta, \quad \alpha + \beta = \tau.
$$

Write

$$
\sigma_U = {}^t\!U\sigma_1 U = \begin{pmatrix} a & b/2 \\ b/2 & c \end{pmatrix}.
$$

Note that σ_U is half-integral, and det $\sigma_U = \det \sigma_1$. So by [\(9.1\)](#page-40-4), either $p \nmid b$ or $p \nmid ac$. We would like to compute the integral

$$
I_{A,p} = \int_{S(\mathbf{Q}_p)} f_p\left(\begin{pmatrix} D & p^\tau SD^{-1} \\ O & p^\tau D^{-1} \end{pmatrix}\right) \theta_p(\text{tr}\,\sigma_U S) dS. \tag{9.5}
$$

Writing $S =$ $\begin{pmatrix} x & y \\ y & z \end{pmatrix}$, we let

$$
M = \begin{pmatrix} D & p^{\tau} S D^{-1} \\ O & p^{\tau} D^{-1} \end{pmatrix} = \begin{pmatrix} p^{\alpha} & 0 & p^{\beta} x & p^{\alpha} y \\ 0 & p^{\beta} & p^{\beta} y & p^{\alpha} z \\ 0 & 0 & p^{\beta} & 0 \\ 0 & 0 & 0 & p^{\alpha} \end{pmatrix}.
$$

By Proposition [3.2,](#page-12-0) $M \in \text{supp } f_p$ if and only if the fractional ideal generated by all the entries is $(1) = \mathbb{Z}_p$ and the fractional ideal generated by the determinants of all 2×2 submatrices is $(p^t) = p^t \mathbb{Z}_p$. The determinants of the 2×2 submatrices of M are shown in the table below:

Using $\alpha \leq \beta$ and $\alpha + \beta = \tau,$ we see that $M \in \operatorname{supp} f_p$ if and only if

$$
(p^{\alpha}, p^{\beta}x, p^{\alpha}y, p^{\alpha}z) = (1)
$$
\n(9.6)

and

$$
(p^{2\alpha}, p^{\tau}x, p^{\tau}y, p^{2\alpha}z, p^{\tau}(xz - y^2)) = (p^t). \tag{9.7}
$$

Let

$$
x' = p^{\beta}x, \quad y' = p^{\alpha}y, \quad z' = p^{\alpha}z.
$$
 (9.8)

Then [\(9.6\)](#page-42-0) is equivalent to

$$
(p^{\alpha}, x', y', z') = (1)
$$
\n(9.9)

and [\(9.7\)](#page-42-1) is equivalent to

$$
(p^{2\alpha}, p^{\alpha}x', p^{\beta}y', p^{\alpha}z', x'z' - p^{\beta - \alpha}y'^2) = (p^t). \tag{9.10}
$$

If $\alpha \neq 0$, then [\(9.6\)](#page-42-0) is equivalent to

$$
(x', y', z') = (1), \tag{9.11}
$$

i.e., $x', y', z' \in \mathbb{Z}_p$ and at least one of them is a unit.

9.2 EVALUATION OF THE INTEGRAL $I_{A,p}$

We continue with the notation from above. Given a Borel subset $S' \subseteq S(\mathbf{Q}_p)$, we define

$$
I_{S'} = \int_{S'} f_p \left(\begin{pmatrix} D & rSD^{-1} \\ O & rD^{-1} \end{pmatrix} \right) \theta_p(\text{tr }\sigma_U S) dS.
$$

Define

$$
S_0 = \left\{ \begin{pmatrix} x & y \\ y & z \end{pmatrix} \in S(\mathbf{Q}_p) \middle| x, y, z \text{ satisfy (9.6) and (9.7)} \right\}.
$$

Then $I_{S_0} = I_{A,p}$ is the integral [\(9.5\)](#page-41-0) we need to compute. Let $E_{11} =$ $\begin{pmatrix} 1 & 0 \\ 0 & 0 \end{pmatrix}$, $E_{22} =$ $\begin{pmatrix} 0 & 0 \\ 0 & 1 \end{pmatrix}$ and E'_{12} = $\begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}$.

PROPOSITION 9.3. Let S' be a Borel subset of S_0 . Suppose $p \nmid a$ (resp. $p \nmid b$, $p \nmid c$). Suppose $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \in S'$ implies that $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \pm \frac{1}{p} E_{11}$ $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \pm \frac{1}{p} E'_{12},$ $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \pm \frac{1}{p} E_{22}$ respectively) belongs to S'. Then $I_{S'} = 0$.

Proof. Suppose $p \nmid a$. The other cases can be handled similarly. By the given property, $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \in S'$ if and only if $\begin{pmatrix} x & y \\ y & z \end{pmatrix} + \frac{1}{p}E_{11} \in S'$. Hence

$$
I_{S'} = \int_{S'} \theta_p(\text{tr }\sigma_U S) dS = \int_{S'} \theta_p(\text{tr }\sigma_U (S - \frac{1}{p}E_{11})) dS = e(\frac{a}{p}) I_{S'}.
$$

The proposition follows.

PROPOSITION 9.4. Suppose $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \in S_0$. Then:

(i) If
$$
\beta \ge 2
$$
, $\tau - 1 \ge t + 1$ and $p^{t+1} | p^{\beta - 1} z'$, then $\begin{pmatrix} x & y \ y & z \end{pmatrix} \pm \frac{1}{p} E_{11} \in S_0$.
\n(ii) If $\alpha \ge 2$, $\tau - 2 \ge t + 1$ and $p^{t+1} | p^{\beta - 1} y'$, then $\begin{pmatrix} x & y \ y & z \end{pmatrix} \pm \frac{1}{p} E'_{12} \in S_0$.
\n(iii) If $\alpha \ge 2$, $2\alpha - 1 \ge t + 1$, $p^{t+1} | p^{\alpha - 1} x'$, then $\begin{pmatrix} x & y \ y & z \end{pmatrix} \pm \frac{1}{p} E_{22} \in S_0$.

Remarks: 1) If in [\(9.10\)](#page-42-2), $(p^{2\alpha}, p^{\alpha}x', p^{\beta}y', p^{\alpha}z') = (p^t)$, then the last condition of (i) (resp. (ii), (iii)) can be replaced by the weaker condition $p^t|p^{\beta-1}z'$ (resp. $p^t|p^{\beta-1}y', p^t|p^{\alpha-1}x'$, and the second condition of (ii) can be weakened to $\tau - 2 \geq t$.

2) There are some other variants; for example, if $(x', z') = 1$ and $\alpha = t$, then in (ii), the conditions can be replaced by $\alpha \geq 1$, $\tau - 2 \geq t$, and $p^t | p^{\beta - 1} y'$.

Proof. (i) Replace x by $x \pm \frac{1}{p}$ in [\(9.6\)](#page-42-0) and [\(9.7\)](#page-42-1). The left-hand side of (9.6) becomes

$$
(p^{\alpha}, p^{\beta}x \pm p^{\beta-1}, p^{\alpha}y, p^{\alpha}z).
$$

The left-hand side of [\(9.7\)](#page-42-1) becomes

$$
(p^{2\alpha}, p^{\tau}x \pm p^{\tau-1}, p^{\tau}y, p^{2\alpha}z, p^{\tau}(xz - y^2) \pm p^{\beta - 1}z').
$$

Under the given hypotheses, $p^{\beta-1} \in (p)$, $p^{\tau-1} \in (p^{t+1})$ and $p^{\beta-1}z' \in (p^{t+1})$. Hence $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \pm \frac{1}{p} E_{11}$ satisfies [\(9.6\)](#page-42-0) and [\(9.7\)](#page-42-1).

(ii) Replace y by $y \pm \frac{1}{p}$ in [\(9.6\)](#page-42-0) and [\(9.7\)](#page-42-1). The left-hand side of (9.6) becomes

$$
(p^{\alpha}, p^{\beta}x, p^{\alpha}y \pm p^{\alpha-1}, p^{\alpha}z).
$$

The left-hand side of [\(9.7\)](#page-42-1) becomes

$$
(p^{2\alpha}, p^{\tau} x, p^{\tau} y \pm p^{\tau-1}, p^{2\alpha} z, p^{\tau} (xz-y^2) \mp 2p^{\beta-1} y' - p^{\tau-2}).
$$

Under the given hypotheses in this case, $p^{\alpha-1} \in (p)$, $p^{\tau-2} \in (p^{t+1})$ and $p^{\beta-1}y' \in$ (p^{t+1}) . Hence $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \pm \frac{1}{p} E'_{12}$ satisfies [\(9.6\)](#page-42-0) and [\(9.7\)](#page-42-1).

Assertion (iii) and the remarks can be proven similarly.

 \Box

COROLLARY 9.5. Suppose $p \nmid a$, and that

- (i) $\beta > 2$, (ii) $\tau - 1 > t + 1$, (iii) $\beta - 1 > t + 1$.
- Then $I_{A,p} = 0$.

Proof. Suppose $\begin{pmatrix} x & y \ y & z \end{pmatrix} \in S_0$. Then by [\(9.9\)](#page-42-3), $z' \in \mathbb{Z}_p$, so by the third hypoth- $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \pm \frac{1}{p} E_{11} \in S_0$. The assertion now esis, $p^{t+1} | p^{\beta - 1} z'$. By Proposition [9.4,](#page-43-0) follows by Proposition [9.3.](#page-43-1) \Box

COROLLARY 9.6. Suppose $p \nmid b$, and

(i) $\alpha \geq 2$, (ii) $\tau - 2 \geq t + 1$, (iii) $\beta - 1 \geq t + 1$.

Then $I_{A,p} = 0$.

Proof. This follows in the same way as the previous corollary, using Proposition [9.4](#page-43-0) (ii), and Proposition [9.3.](#page-43-1) \Box

PROPOSITION 9.7. Suppose condition (iii) of the above corollaries fails to hold, i.e., $\beta - 1 \leq t$. Then exactly one of the following is true:

1. $\tau = 2\tau' + 1$ is odd, $\alpha = \tau'$, $\beta = \tau' + 1$, and $t = \tau'$, 2. $\tau = 2\tau'$ is even, $\alpha = \tau' - 1$, $\beta = \tau' + 1$, and $t = \tau'$, 3. $\tau = 2\tau'$ is even, and $\alpha = \beta = t = \tau'$, 4. $\tau = 2\tau'$ is even, $\alpha = \beta = \tau'$, and $t = \tau' - 1$.

Proof. Suppose $\beta \leq t + 1$. Then because we always have $t \leq \left[\frac{\tau}{2}\right]$ (where brackets denote the floor), it follows that $\beta \leq t+1 \leq \left[\frac{\tau}{2}\right]+1$. On the other hand, $\tau = \alpha + \beta \leq 2\beta$, which gives the lower bound in the following inequality:

$$
\left\lceil\frac{\tau}{2}\right\rceil\leq\beta\leq\left\lceil\frac{\tau}{2}\right\rceil+1.
$$

(Here, $\lceil \cdot \rceil$ denotes the ceiling.) Using $\beta - 1 \le t \le \lceil \frac{\tau}{2} \rceil$, the result follows easily by considering the possible cases.

PROPOSITION 9.8. Suppose α, β, t, τ satisfy Proposition [9.7](#page-44-0) (1), i.e., $\tau = 2\tau' +$ 1, $\alpha = \tau'$, $\beta = \tau' + 1$, and $t = \tau'$. Then if $\tau' \geq 2$, $I_{A,p} = 0$.

Proof. Let $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \in S_0$. Then [\(9.10\)](#page-42-2) is satisfied, and since $\alpha = \tau' \ge 2$, [\(9.11\)](#page-42-4) is also satisfied. In particular, by [\(9.10\)](#page-42-2),

$$
x'z' \equiv py'^2 \pmod{p^2}.
$$
\n(9.12)

It follows that either x' or z' is a unit. Indeed, if $p|x'$ and $p|z'$, then y' is a unit by (9.11) , leading to an obvious contradiction in (9.12) . In fact, by (9.12) , $p|x'z'$ and hence exactly one of x' or z' is a unit.

First suppose $p \nmid b$. Note that $(x', z') = 1$, $\alpha = t \geq 2$,

$$
\tau - 2 = 2\tau' - 1 > 2\tau' - \tau' = t,
$$

and $p^t|p^{\beta-1}y'$. By the second remark after Proposition [9.4,](#page-43-0) $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \pm \frac{1}{p} E'_{12} \in$ S_0 . By Proposition [9.3,](#page-43-1) $I_{A,p} = I_{S_0} = 0$.

Finally, suppose p|b. Then as noted earlier, $p \nmid a$. Since one of x' or z' is a unit, $(p^{\alpha}x', p^{\alpha}z') = (p^t)$. Furthermore, $p^t|p^{\beta-1}z'$, and as above, $\tau - 2 \ge t$. By the $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \pm \frac{1}{p} E_{11} \in S_0$. By Proposition [9.3,](#page-43-1) first remark after Proposition [9.4,](#page-43-0) $I_{A,p} = I_{S_0} = 0.$ \Box

PROPOSITION 9.9. Suppose α, β, t, τ satisfy Proposition [9.7](#page-44-0) (2), i.e., $\tau = 2\tau'$ is even, $\alpha = \tau' - 1$, $\beta = \tau' + 1$, and $t = \tau'$. Then if $\tau' \geq 3$, $I_{A,p} = 0$. In fact, if $p \nmid ac$, then $I_{A,p} = 0$ if $\tau' \geq 2$.

Proof. Let $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \in S_0$. Suppose $\tau' \geq 2$. By [\(9.10\)](#page-42-2), $p^{\tau'}|p^{\tau'-1}x'$ and $p^{\tau'}|p^{\tau'-1}z'$, and hence $p|x'$ and $p|z'$. Therefore by [\(9.11\)](#page-42-4), y' is a unit.

Suppose $p \nmid ac$. Because $p|z'$, we have $p^{t+1}|p^{\beta-1}z'$. Hence by Proposition [9.4,](#page-43-0) $\left(\begin{array}{c} x \ y \ z\end{array}\right) \pm \frac{1}{p}E_{11} \in S_0$. By Proposition [9.3,](#page-43-1) $I_{S_0} = 0$.

Next suppose $p \nmid b$ and $\tau' \geq 3$. Write $x' = px''$ and $z' = pz''$, with x'' and $z'' \in \mathbb{R}$ \mathbf{Z}_p . By [\(9.10\)](#page-42-2), $p^2 x'' z'' \equiv p^2 y'^2 \pmod{p^t}$, so $x'' z'' \equiv y'^2 \pmod{p^{\tau'-2}}$. Because

 $\tau' \geq 3$, it follows that x'' and z'' are units. Therefore $(p^{\alpha}x') = (p^t)$. Obviously $p^t|p^{\beta-1}y'$. By the first remark after Proposition [9.4,](#page-43-0) $\begin{bmatrix} x & y \\ y & z \end{bmatrix} \pm \frac{1}{p}E'_{12} \in S_0$. By Proposition [9.3,](#page-43-1) $I_{A,p} = I_{S_0} = 0$.

PROPOSITION 9.10. Suppose α, β, t, τ satisfy Proposition [9.7](#page-44-0) (3), i.e. $\tau = 2\tau'$ and $\alpha = \beta = t = \tau'$. Suppose further that $\tau' \geq 2$. Then the integral $I_{A,p}$ is given explicitly by [\(9.14\)](#page-47-0) below.

Proof. Suppose $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \in S_0$. Then [\(9.11\)](#page-42-4) implies $(p^t x', p^t y', p^t z') = (p^t)$. Hence (9.11) and (9.10) taken together are equivalent to (9.11) and

$$
x'z' \equiv y'^2 \pmod{p^t}.\tag{9.13}
$$

If y' is not a unit, then by (9.13) and (9.11) , exactly one of x' or z' is a unit. So there is a partition

$$
S_0 = S_1 \cup S_2 \cup S_3,
$$

where

$$
S_1 = \left\{ \begin{pmatrix} x & y \\ y & z \end{pmatrix} \in S_0 : x', y', z' \text{ are units} \right\},
$$

\n
$$
S_2 = \left\{ \begin{pmatrix} x & y \\ y & z \end{pmatrix} \in S_0 : p|y', p|z', \text{ and } x' \text{ is a unit} \right\},
$$

\n
$$
S_3 = \left\{ \begin{pmatrix} x & y \\ y & z \end{pmatrix} \in S_0 : p|y', p|x', \text{ and } z' \text{ is a unit} \right\}.
$$

We claim that $I_{S_2} = I_{S_3} = 0$. Let $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \in S_2$. Then $(p^{\alpha}x') = (p^t), p^t|p^{\beta-1}z'$, and $\tau-1 \geq t+1$. By the first remark after Proposition [9.4,](#page-43-0) $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \pm \frac{1}{p} E_{11} \in S_0.$ In fact, this matrix belongs to S_2 since $p^{\beta}(x \pm \frac{1}{p}) = x' \pm p^{\beta - 1}$ is a unit. Hence by Proposition [9.3,](#page-43-1) $I_{S_2} = 0$ if $p \nmid ac$. If $p \mid ac$, then $p \nmid b$. In this case, $(p^{\alpha}x') = (p^t)$, $p^t \lbrack p^{\beta-1}y'$, and $\tau - 1 \geq t+1$. By the first remark after Proposition [9.4,](#page-43-0) $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \pm \frac{1}{p} E'_{12} \in S_0$. In fact, the matrix belongs to S_2 since $p^{\alpha}(y \pm \frac{1}{p}) = y' \pm p^{\alpha-1} \in p\mathbb{Z}_p$. By Proposition [9.3,](#page-43-1) $I_{S_2} = 0$. The proof that $I_{S_3} = 0$ is similar.

For the integral over S_1 , note that the validity of (9.11) and (9.13) depends only on $x', y', z' \pmod{p^{\tau'}}$, which also means that $S_1 = S_1 + S_2(\mathbf{Z}_p)$ (where S_2 here denotes the symmetric matrices). Hence, writing x, y, z for the congruence classes of $x', y', z' \mod p^{\tau'}$,

$$
I_{A,p} = I_{S_1} = \int_{S_1} \theta_p(\text{tr}(\sigma_U S))dS
$$

$$
\sum_{\substack{x,y,z \in (\mathbf{Z}_p/p^{\tau'}\mathbf{Z}_p)^* , \int_{\mathbf{Z}^{\tau'}\mathbf{Z}_p}} \int_{\frac{z}{p^{\tau'}} + \mathbf{Z}_p} \int_{\frac{y}{p^{\tau'}} + \mathbf{Z}_p} \int_{\frac{x}{p^{\tau'}} + \mathbf{Z}_p} \theta_p(ax + by + cz) dx dy dz
$$
\n
$$
= \sum_{\substack{y \in (\mathbf{Z}_p/p^{\tau'}\mathbf{Z}_p)^* \\ y \in (\mathbf{Z}_p/p^{\tau'}\mathbf{Z}_p)^*}} e(-\frac{by}{p^{\tau'}}) \left(\sum_{\substack{x,z \in \mathbf{Z}_p/p^{\tau'}\mathbf{Z}_p \\ x \equiv y^2 (\text{mod } p^{\tau'})}} e(-\frac{ax + cz}{p^{\tau'}}) \right)
$$
\n
$$
= \sum_{\substack{y \in (\mathbf{Z}_p/p^{\tau'}\mathbf{Z}_p)^* \\ x \equiv 2 (\text{mod } p^{\tau'})}} e(-\frac{by}{p^{\tau'}}) \left(\sum_{\substack{x,z \in \mathbf{Z}_p/p^{\tau'}\mathbf{Z}_p \\ x \equiv 1 (\text{mod } p^{\tau'})}} e(-\frac{y(ax + cz)}{p^{\tau'}}) \right)
$$
\n
$$
= \sum_{\substack{x,z \in \mathbf{Z}_p/p^{\tau'}\mathbf{Z}_p \\ x \equiv 1 (\text{mod } p^{\tau'})}} \left(\sum_{y \in (\mathbf{Z}_p/p^{\tau'}\mathbf{Z}_p)^*} e(-\frac{y(ax + cz + b)}{p^{\tau'}}) \right).
$$
\n(9.14)

We remark that the computation here is valid for $\tau' \geq 1$. The sum over y can be evaluated using [\(9.20\)](#page-48-0) below.

PROPOSITION 9.11. Suppose α, β, t, τ satisfy Proposition [9.7](#page-44-0) (4), i.e. $\tau = 2\tau'$ is even, $\alpha = \beta = \tau'$, and $t = \tau' - 1$. Suppose further that $\tau' \geq 2$. Then the integral $I_{A,p}$ is given by [\(9.18\)](#page-48-1) below.

Proof. In this case, (9.11) and (9.10) are equivalent to (9.11) and

$$
(x'z'-{y'}^2)=(p^{\tau'-1}),
$$
\n(9.15)

i.e.,

$$
x'z' \equiv y'^2 \pmod{p^{\tau'-1}}
$$
\n
$$
(9.16)
$$

but

$$
x'z' \not\equiv y'^2 \pmod{p^{\tau'}}
$$
\n
$$
(9.17)
$$

As in the previous proof, we integrate over S_1, S_2 and S_3 . We first show that $I_{S_2} = I_{S_3} = 0$. Suppose $\left(\begin{matrix} x & y \\ y & z \end{matrix}\right) \in S_2$, so $x' \in \mathbb{Z}_p^*$, $p|y'$, and $p|z'$. Then the conditions of Proposition [9.4\(](#page-43-0)i) are satisfied, so $\begin{pmatrix} x & y \\ y & z \end{pmatrix} \pm \frac{1}{p}E_{11} \in S_0$. Since $p^{\beta}(x+\frac{1}{p})=x'+p^{\beta-1}$ is a unit, this matrix in fact belongs to S_2 . By Proposition [9.3,](#page-43-1) $I_{S_2} = 0$, assuming $p \nmid ac$. On the other hand, if $p \mid ac$, then $p \nmid b$. The hypotheses to Proposition [9.4\(](#page-43-0)ii) are satisfied, so $\begin{pmatrix} x & y \ y & z \end{pmatrix} \pm \frac{1}{p} E'_{12} \in S_2$, and once again Proposition [9.3](#page-43-1) gives $I_{S_2} = 0$. The proof that $I_{S_3} = 0$ is similar. For the integral over S_1 , just as in the proof of the previous proposition, we have

$$
I_{A,p} = I_{S_1} = \int_{S_1} \theta_p(\text{tr}\,\sigma_U S) dS
$$

Documenta Mathematica 24 (2019) 677–747

=

$$
= \sum_{\substack{x,y,z \in (\mathbf{Z}_{p}/p^{\tau'}\mathbf{Z}_{p})^{*},\\ x \equiv y^{2}(\text{mod } p^{\tau'-1}),\ xz \not\equiv y^{2}(\text{mod } p^{\tau'})}} \int_{\frac{z}{p^{\tau'}} + \mathbf{Z}_{p}} \int_{\frac{y}{p^{\tau'}} + \mathbf{Z}_{p}} \int_{\frac{x}{p^{\tau'}} + \mathbf{Z}_{p}} \theta_{p}(ax + by + cz) dxdydz
$$

\n
$$
= \sum_{y \in (\mathbf{Z}_{p}/p^{\tau'}\mathbf{Z}_{p})^{*}} e(-\frac{by}{p^{\tau'}}) \left(\sum_{\substack{x,z \in \mathbf{Z}_{p}/p^{\tau'}\mathbf{Z}_{p} \\ x \equiv y^{2}(\text{mod } p^{\tau'-1}),\ xz \not\equiv y^{2}(\text{mod } p^{\tau'})}} e(-\frac{ax + cz}{p^{\tau'}}) \right)
$$

\n
$$
= \sum_{y \in (\mathbf{Z}_{p}/p^{\tau'}\mathbf{Z}_{p})^{*}} e(-\frac{by}{p^{\tau'}}) \left(\sum_{\substack{x,z \in \mathbf{Z}_{p}/p^{\tau'}\mathbf{Z}_{p} \\ xz \equiv 1(\text{mod } p^{\tau'-1}),\ xz \not\equiv 1(\text{mod } p^{\tau'}-1),\ xz \not\equiv 1(\text{mod } p^{\tau'})}} e(-\frac{y(ax + cz)}{p^{\tau'}}) \right)
$$

\n
$$
= \sum_{\substack{x,z \in \mathbf{Z}_{p}/p^{\tau'}\mathbf{Z}_{p} \\ xz \equiv 1(\text{mod } p^{\tau'-1}),\ xz \not\equiv 1(\text{mod } p^{\tau'})}} \left(\sum_{y \in (\mathbf{Z}_{p}/p^{\tau'}\mathbf{Z}_{p})^{*}} e(-\frac{y(ax + cz + b)}{p^{\tau'}}) \right) (9.18)
$$

\n
$$
= \sum_{h=1}^{p-1} \sum_{\substack{x,z \in \mathbf{Z}_{p}/p^{\tau'}\mathbf{Z}_{p} \\ xz \equiv 1 + h p^{\tau'-1}(\text{mod } p^{\tau'})}} \left(\sum_{y \in (\mathbf{Z}_{p}/p^{\tau'}\mathbf{Z}_{p})^{*}} e(-\frac{y(ax + cz + b)}{p^{\tau'}}) \right).
$$

Once again, this computation is valid for $\tau' \geq 1$. We remark that [\(9.18\)](#page-48-1) can be rewritten as

$$
\sum_{\substack{x,z \in \mathbf{Z}_p/p^{\tau'}\mathbf{Z}_p \\ xz \equiv 1 (\text{mod } p^{\tau'-1})}} \left(\sum_{\mathbf{y} \in (\mathbf{Z}_p/p^{\tau'}\mathbf{Z}_p)^*} e\left(-\frac{\mathbf{y}(a\mathbf{x} + c\mathbf{z} + b)}{p^{\tau'}}\right) - (9.14). \tag{9.19}
$$

COROLLARY 9.12. Suppose $\beta - 1 \leq t$ and $\tau \geq 5$. Then $I_{A,p} \ll_p p^{3\tau/4}$.

Proof. By Propositions [9.7-](#page-44-0)[9.11,](#page-47-1) we may assume that $\tau = 2\tau'$ is even, and $I_{A,p}$ is given by [\(9.14\)](#page-47-0) or [\(9.18\)](#page-48-1). Recall the formula for the Ramanujan sum

$$
\sum_{\mathbf{y}\in(\mathbf{Z}_p/p^{\tau'}\mathbf{Z}_p)^*} e(\frac{\mathbf{y}\ell}{p^{\tau'}}) = \begin{cases} p^{\tau'} - p^{\tau'-1} & \text{if } p^{\tau'}|\ell, \\ -p^{\tau'-1} & \text{if } p^{\tau'-1}||\ell, \\ 0 & \text{if } p^{\tau'-1} \nmid \ell \end{cases}
$$
(9.20)

([\[Hua,](#page-68-1) Theorem 4.3]). Thus the summation over y in [\(9.14\)](#page-47-0) is nonzero only if

$$
ax + cz + b \equiv 0 \pmod{p^{\tau'-1}}.
$$
\n(9.21)

Since $xz \equiv 1 \mod p^{\tau'}$, this is equivalent to

$$
ax^{2} + bx + c \equiv 0 \mod p^{\tau'-1},
$$
 (9.22)

and also to

$$
cz^2 + bz + a \equiv 0 \mod p^{\tau'-1}.
$$
 (9.23)

Suppose $p|a$ and $p|c$, so that $p \nmid b$. Then [\(9.21\)](#page-48-2) has no solution, so the summation is zero. If $p \nmid a$ (resp. $p \nmid c$), then the number of x mod $p^{\tau'}$ satisfying [\(9.22\)](#page-48-3) (resp. z mod $p^{\tau'}$ satisfying [\(9.23\)](#page-49-1)) is $\ll pp^{(\tau'-1)/2} \ll_p p^{\tau'/2}$ (cf. [\[KL3,](#page-68-2) Lemma 9.6]). Now applying [\(9.20\)](#page-48-0), we see that [\(9.14\)](#page-47-0) is $O(p^{3\tau'/2})$. Similarly, [\(9.18\)](#page-48-1) is

$$
\ll \sum_{h=1}^{p-1} \sum_{\substack{x,z \in \mathbb{Z}_p/p^{\tau'} \mathbb{Z}_p \\ x \equiv 1 + h p^{\tau'-1} \pmod{p^{\tau'}}, \\ ax + cz + b \equiv 0 \pmod{p^{\tau'-1}}} p^{\tau'} \ll p^2 p^{\frac{\tau'-1}{2}} p^{\tau'} \ll p^{3\tau'/2}.
$$

9.3 PROOF OF PROPOSITION [9.1](#page-40-3)

For $p \in S$, we require a bound for the local integral $I_{A,p}$ given in [\(9.5\)](#page-41-0). We continue to use the notation of Section [9.1.](#page-41-1) Thus, $\tau = r_p$, $\lambda(p)$ = $\text{diag}(1, p^t, p^{\tau}, p^{\tau-t})$ for $0 \leq t \leq \tau/2$, $D = \text{diag}(p^{\alpha}, p^{\beta})$ for $0 \leq \alpha \leq \beta$ with $\alpha + \beta = \tau$, and $\sigma_U = {}^t\!U\sigma_1 U = \begin{pmatrix} a & b/2 \\ b/2 & a \end{pmatrix}$ $b/2$ c \setminus . By [\(3.10\)](#page-9-0)

$$
\rho = \frac{3}{2}e_0 - 2e_1 - e_2,
$$

so

$$
\langle \lambda, \rho \rangle = \frac{3}{2} \ell_0 - 2 \ell_1 - \ell_2 = \frac{3}{2} \tau - t.
$$

To prove Proposition [9.1,](#page-40-3) we must show that for some $\varepsilon > 0$,

$$
|I_{A,p}(c_{\lambda})| \ll p^{(1-\varepsilon)(\lambda,\rho)-\varepsilon\tau} = p^{(1-\varepsilon)(\frac{3}{2}\tau-t)-\varepsilon\tau}.
$$
 (9.24)

By the trivial bound [\(7.13\)](#page-32-0),

$$
|I_{A,p}|\leq p^{2\alpha+\beta}=p^{\alpha+\tau}.
$$

Therefore [\(9.24\)](#page-49-2) is certainly satisfied when

$$
\alpha + \tau \le (1 - \varepsilon)(\frac{3}{2}\tau - t) - \varepsilon\tau,
$$
\n(9.25)

or equivalently,

$$
\alpha + (1 - \varepsilon)t \le \frac{1 - 5\varepsilon}{2}\tau.
$$

Note that $(1 - \varepsilon)(\frac{3}{2}\tau - t) - \varepsilon\tau$ is a decreasing function of ε . Therefore if (9.24) holds for some particular $\varepsilon = \varepsilon_0$, then it holds for all smaller positive ε . For

concreteness, we will verify it for $\varepsilon_0 = 0.01$, in which case the above inequality takes the form

$$
\alpha + 0.99t \le 0.475\tau. \tag{9.26}
$$

For any given value of τ , there are only finitely many permissible values for t, α and β , so the associated integral is bounded by a constant depending only on τ . Therefore we may assume that

$$
\tau \ge 7. \tag{9.27}
$$

.

If $p \nmid b$ (resp. $p \nmid ac$) and the conditions of Corollary [9.5](#page-44-1) (resp. Corollary [9.6\)](#page-44-2) hold, the integral vanishes and the desired bound is trivially satisfied.

Suppose condition (iii) of either Corollary [9.5](#page-44-1) or Corollary [9.6](#page-44-2) fails. Then by Corollary [9.12,](#page-48-4)

$$
I_{A,p} \ll p^{3\tau/4} < p^{(1-2\varepsilon)\tau}
$$

when $\varepsilon < 1/8$. If $t = \tau/2$ (as is the case in Proposition [9.10\)](#page-46-1), then

$$
p^{(1-2\varepsilon)\tau} = p^{(1-\varepsilon)(\frac{3}{2}\tau - t) - \varepsilon\tau}
$$

If $t = \tau/2 - 1$ (as is the case in Proposition [9.11\)](#page-47-1), then

$$
p^{(1-2\varepsilon)\tau} \ll p^{(1-2\varepsilon)\tau + (1-\varepsilon)} = p^{(1-\varepsilon)(\frac{3}{2}\tau - t) - \varepsilon\tau}.
$$

Either way, we obtain the desired bound for $I_{A,p}$ when $\varepsilon < 1/8$. Suppose condition (i) of either Corollary [9.5](#page-44-1) or Corollary [9.6](#page-44-2) fails. Then (using $\alpha \leq \beta$ in the first case) $\alpha \leq 1$. By [\(9.7\)](#page-42-1), $t \leq 2$. Hence by [\(9.27\)](#page-50-0),

 $\alpha + 0.99t < 3 < 3.325 = 0.475 \times 7 \le 0.475\tau$,

so [\(9.26\)](#page-50-1) is satisfied in this case, and the desired bound holds. Suppose condition (ii) of either Corollary [9.5](#page-44-1) or Corollary [9.6](#page-44-2) fails. Then $\tau - 2 \le t \le \frac{\tau}{2}$, which means that $\tau \le 4$, contradicting [\(9.27\)](#page-50-0). This proves [\(9.24\)](#page-49-2) and hence Proposition [9.1.](#page-40-3)

A DISCRETE SERIES MATRIX COEFFICIENTS FOR $GSp(2n)$

Here we explicitly compute certain discrete series matrix coefficients for $GSp_{2n}(\mathbf{R})$ using ideas of Harish Chandra. Our main references for the background material are [\[AS\]](#page-67-3) and [\[Kn\]](#page-68-3).

A.1 Root System

The Lie algebra of $\text{Sp}_{2n}(\mathbf{R})$ is

$$
\mathfrak{g} = \{X | JX + {}^tXJ = 0\}
$$

=
$$
\left\{ \begin{pmatrix} A & B \\ C & D \end{pmatrix} \in M_{2n}(\mathbf{R}) | A = -{}^tD, B = {}^tB, C = {}^tC \right\}.
$$

We have

$$
\mathfrak{g}=\mathfrak{k}\oplus\mathfrak{p},
$$

where

$$
\mathfrak{k} = \left\{ \begin{pmatrix} A & B \\ -B & A \end{pmatrix} \in M_{2n}(\mathbf{R}) | A = -{}^{t}A, B = {}^{t}B \right\}
$$

is the Lie algebra of $K,^{\dagger}$

$$
\mathfrak{p} = \left\{ \begin{pmatrix} A & B \\ B & -A \end{pmatrix} \in M_{2n}(\mathbf{R}) \middle| A = {}^{t}A, B = {}^{t}B \right\}.
$$

Let $\mathfrak{h} \subseteq \mathfrak{k}_{\mathbb{C}}$ denote the real subspace consisting of all matrices of the form

 a¹ a² . . . ^aⁿ [−]a¹ −a² . . . −aⁿ , (A.1)

for $a_j \in i\mathbf{R}$ ($i^2 = -1$). Then h is a compact Cartan subalgebra of $\mathfrak{g}_\mathbf{C}$. For each $j = 1, \ldots, n$, define a linear form e_j on \mathfrak{h}_c by taking the above matrix to ia_j . Let $\Delta = \Delta(\mathfrak{h}_{\mathbf{C}}, \mathfrak{g}_{\mathbf{C}}) \subseteq \mathfrak{h}_{\mathbf{C}}^*$ be the set of roots. One finds that

$$
\Delta = \{ \pm 2e_j, \pm (e_j + e_k), e_j - e_k | j \neq k \}.
$$

The compact roots (i.e. those whose root spaces belong to $\mathfrak{k}_{\mathbb{C}}$) are

$$
\Delta_K = \{e_j - e_k | j \neq k\}.
$$

Let $\Delta_{nc} = \Delta - \Delta_K$ denote the set of noncompact roots. We fix the inner product in Δ determined by

$$
\langle e_i, e_j \rangle = \delta_{ij}.
$$

We set $\varepsilon_j = -e_j$, and fix the following ordered basis for $\mathfrak{h}_{\mathbf{C}}^*$:

 $\{\varepsilon_1+\varepsilon_2+\cdots+\varepsilon_n, \varepsilon_n-\varepsilon_{n-1}, \varepsilon_{n-1}-\varepsilon_{n-2}, \ldots, \varepsilon_2-\varepsilon_1\}.$

This determines a good ordering of Δ (i.e. $\Delta_{nc}^+ > \Delta_K^+$) in which $\Delta^+ = \Delta_{nc}^+ \cup \Delta_K^+$ is given by

$$
2\varepsilon_j, \quad 1 \le j \le n
$$

\n
$$
\varepsilon_j + \varepsilon_k, \quad 1 \le j < k \le n
$$

\n
$$
\varepsilon_k - \varepsilon_j, \quad 1 \le j < k \le n
$$

\n
$$
\Delta_K^+.
$$

[†]In this appendix, K denotes the compact subgroup [\(3.4\)](#page-8-0) which was denoted K_{∞} elsewhere in the paper.

Then

$$
\delta_G = \frac{1}{2} \sum_{\alpha \in \Delta^+} \alpha = \varepsilon_1 + 2\varepsilon_2 + \ldots + n\varepsilon_n.
$$
 (A.2)

Remark: The positive system (here denoted $\Delta^{'}{}^+$) considered in [\[AS\]](#page-67-3) is the one used in Chapter IX of [\[Kn\]](#page-68-3). The relationship to the present system is $\Delta'_{K}^+=\Delta^+_{K}$ and $\Delta'_{nc}^+=-\Delta^+_{nc}$. In other words, $\Delta^+=-w_{K}(\Delta'^{+})$, where w_{K} is the element of the Weyl group W_K taking Δ_K^+ to $-\Delta_K^+$. (See Remark (1) after Theorem 9.20 in [\[Kn\]](#page-68-3).)

For each root $\alpha \in \Delta$, let $\mathfrak{g}_{\alpha} \subseteq \mathfrak{g}_{\mathbb{C}}$ be the associated root space. Define the following subalgebras of $\mathfrak{g}_{\mathbb{C}}$:

$$
\mathfrak{p}^+ = \bigoplus_{\alpha \in \Delta_{nc}^+} \mathfrak{g}_{\alpha} = \left\{ \begin{pmatrix} A & -iA \\ -iA & -A \end{pmatrix} | A = {}^{t}A \in M_n(\mathbf{C}) \right\},\,
$$

and

$$
\mathfrak{p}^- = \bigoplus_{\alpha \in \Delta_{nc}^+} \mathfrak{g}_{-\alpha} = \left\{ \begin{pmatrix} A & iA \\ iA & -A \end{pmatrix} | A = {}^t A \in M_n(\mathbf{C}) \right\}.
$$

The corresponding analytic subgroups of $Sp_{2n}(\mathbb{C})$ are

$$
P^{+} = \exp(\mathfrak{p}^{+}) = \left\{ \begin{pmatrix} I_{n} + A & -iA \\ -iA & I_{n} - A \end{pmatrix} | A = {}^{t}A \in M_{n}(\mathbf{C}) \right\}
$$

and

$$
P^{-} = \exp(\mathfrak{p}^{-}) = \left\{ \begin{pmatrix} I_n + A & iA \\ iA & I_n - A \end{pmatrix} | A = {}^{t}A \in M_n(\mathbf{C}) \right\}.
$$

We also have

$$
K_{\mathbf{C}} = \exp(\mathfrak{k}_{\mathbf{C}}) = \left\{ \begin{pmatrix} A & B \\ -B & A \end{pmatrix} \in GL_{2n}(\mathbf{C}) | (A + iB) \; \mathfrak{k}(A - iB) = I_n \right\}.
$$

Since A and B are complex, the condition on $(A+iB)$ is equivalent to $A+iB \in$ $GL_n(\mathbf{C})$, reflecting the fact that $U(n, \mathbf{C}) \cong GL_n(\mathbf{C})$.

A.2 REALIZATION IN $SU(n, n)$

Recall that

$$
SU(n,n) = \left\{ g \in SL_{2n}(\mathbf{C}) \Big| \, {}^t \overline{g} \begin{pmatrix} I_n & \\ & -I_n \end{pmatrix} g = \begin{pmatrix} I_n & \\ & -I_n \end{pmatrix} \right\}.
$$

Define

$$
G' = \{ g \in \mathrm{SU}(n,n) | \,^t g J g = J \}.
$$

One can show that

$$
G' = \left\{ \left(\frac{\alpha}{\beta} \quad \frac{\beta}{\alpha} \right) \in SL_{2n}(\mathbf{C}) \middle| \quad \substack{t \overline{\alpha}\alpha - \, t \beta \overline{\beta} = I_n \\ t \overline{\beta}\alpha = \, t \alpha \overline{\beta}} \right\}.
$$

Let $\tau = \begin{pmatrix} I_n & iI_n \\ iI_n & I_n \end{pmatrix}$. Then the map

$$
g \mapsto g' = \tau^{-1} g \tau
$$

is an isomorphism from $\text{Sp}_{2n}(\mathbf{R})$ into G'. For any object O associated to $\text{Sp}_{2n}(\mathbf{R})$, we let O' denote the corresponding object for G'. Writing $g = (A \ B)$ $\begin{pmatrix} A & B \\ C & D \end{pmatrix}$, we have

$$
g' = \frac{1}{2} \begin{pmatrix} (A+D) + i(B-C) & (B+C) + i(A-D) \\ (B+C) - i(A-D) & (A+D) - i(B-C) \end{pmatrix}.
$$
 (A.3)

Taking $g =$ $\begin{pmatrix} A & B \\ -B & A \end{pmatrix} \in K$, we see that $g' = \begin{pmatrix} A+Bi & A-Bi \end{pmatrix}$, where $A + Bi$ is unitary. Thus

$$
K' = \tau^{-1} K \tau = \left\{ \begin{pmatrix} \alpha & \mathbf{1}_{\alpha^{-1}} \end{pmatrix} \mid \alpha \in \mathcal{U}(n) \right\}.
$$

Likewise,

$$
K'_{\mathbf{C}} = \tau^{-1} K_{\mathbf{C}} \tau = \left\{ \begin{pmatrix} \alpha & \mathbf{t}_{\alpha^{-1}} \end{pmatrix} \mid \alpha \in \text{GL}_n(\mathbf{C}) \right\}
$$

$$
P'^{+} = \tau^{-1} P^{+} \tau = \left\{ \begin{pmatrix} I_n & O \ -2iA & I_n \end{pmatrix} \mid A = \,{}^{t}A \in M_n(\mathbf{C}) \right\}
$$

and

$$
P'^{-} = \tau^{-1} P^{-} \tau = \left\{ \begin{pmatrix} I_n & 2iA \\ O & I_n \end{pmatrix} | A = {}^{t}A \in M_n(\mathbf{C}) \right\}.
$$

A.3 HOLOMORPHIC DISCRETE SERIES

We recall without proof some properties of the holomorphic discrete series for the group $G = \mathrm{Sp}_{2n}(\mathbf{R})$. This material is due to Harish-Chandra [\[HC\]](#page-67-4). We follow the exposition in Chapter VI of [\[Kn\]](#page-68-3). Suppose $\lambda \in \mathfrak{h}_{\mathbf{C}}^*$ is analytically integral (i.e. $\lambda(H) \in 2\pi i \mathbb{Z}$ whenever $\exp(H) = 1$) and dominant with respect to K (i.e. $\langle \lambda, \alpha \rangle > 0$ for all $\alpha \in \Delta_K^+$). Let (Φ_λ, V) denote the irreducible unitary representation of K with highest weight λ . Let $v_{\lambda} \in V$ be a highest weight unit vector (unique up to unitary scaling).

Extend Φ_{λ} to a holomorphic representation of $K_{\mathbf{C}}$. For $g \in G$, let $\mu(g) \in K_{\mathbf{C}}$ denote the middle component in the Harish-Chandra decomposition

$$
G \subseteq P^+ K_{\mathbf{C}} P^-
$$

([\[Kn,](#page-68-3) Theorem 6.3]). Define

$$
\psi_{\lambda}(g) = \left\langle \Phi_{\lambda}(\mu(g))^{-1} v_{\lambda}, v_{\lambda} \right\rangle_{V}.
$$

Then under the condition

$$
\langle \lambda + \delta_G, \alpha \rangle < 0 \quad \text{ for all } \quad \alpha \in \Delta_{nc}^+, \tag{A.4}
$$

 ψ_{λ} is a nontrivial square-integrable function on G ([\[Kn,](#page-68-3) Lemma 6.9]), and its translates under the left regular representation generate an irreducible squareintegrable representation $(\pi_{\lambda}, V_{\lambda})$ of G ([\[Kn,](#page-68-3) Theorem 6.6]). Furthermore, by [\[Kn,](#page-68-3) p. 160 (6)],

$$
\langle \pi_{\lambda}(g)\psi_{\lambda}, \psi_{\lambda}\rangle_{L^{2}} = \psi_{\lambda}(g^{-1})\|\psi_{\lambda}\|^{2}.
$$

To each $v \in V$ we associate the function $\langle \Phi_\lambda(\mu(g))^{-1}v, v_\lambda \rangle \in V_\lambda$. This defines a K-equivariant embedding $V \to V_{\lambda}$, and hence Φ_{λ} occurs as a K-type in π_{λ} , with highest weight vector ψ_{λ} . In fact, this K-type occurs with multiplicity oneand $w_{\lambda} = \frac{\psi_{\lambda}}{\|\psi_{\lambda}\|} \in V_{\lambda}$ is a highest weight unit vector ([\[Kn,](#page-68-3) p. 160 (5)]). Our aim is to compute the matrix coefficient

$$
\langle \pi_{\lambda}(g)w_{\lambda}, w_{\lambda} \rangle = \psi_{\lambda}(g^{-1}) = \langle \Phi_{\lambda}(\mu(g^{-1})^{-1})v_{\lambda}, v_{\lambda} \rangle_{V}. \tag{A.5}
$$

We will use the realization of G in $SU(n, n)$ since it facilitates the computation of $\mu(g)$. Let $g' = \begin{pmatrix} \alpha & \beta \\ \frac{\alpha}{2} & \frac{\beta}{2} \end{pmatrix}$ β $\overline{\alpha}$ $\Big(\big) \in G'.$ Then the Harish-Chandra decomposition is given explicitly by

$$
\begin{pmatrix}\n\alpha & \beta \\
\overline{\beta} & \overline{\alpha}\n\end{pmatrix} =\n\begin{pmatrix}\nI_n & O \\
\overline{\beta}\alpha^{-1} & I_n\n\end{pmatrix}\n\begin{pmatrix}\n\alpha & O \\
O & t_{\alpha^{-1}}\n\end{pmatrix}\n\begin{pmatrix}\nI_n & \alpha^{-1}\beta \\
O & I_n\n\end{pmatrix}.
$$
\n(A.6)

To verify this decomposition, note that the lower right corner on the right-hand side is ${}^t\!\alpha^{-1} + \overline{\beta}\alpha^{-1}\beta$. By the fact that ${}^t\!\alpha\overline{\alpha} = I_n + \overline{{}^t\!\beta}\beta$,

$$
\overline{\alpha} = {}^{t} \alpha^{-1} + {}^{t} \alpha^{-1} \overline{{}^{t} \beta} \beta.
$$

Also $\overline{t\beta}\alpha = t_{\alpha\beta} \implies t_{\alpha^{-1}} = \overline{\beta}\alpha^{-1}\overline{t_{\beta}}^{-1}$. Substituting this into the second term above, we see that the lower right-hand corner is equal to $\overline{\alpha}$ as needed.

Let
$$
g = \begin{pmatrix} A & B \\ C & D \end{pmatrix} \in \text{Sp}_{2n}(\mathbf{R})
$$
, and let $g' = \begin{pmatrix} \alpha & \beta \\ \overline{\beta} & \overline{\alpha} \end{pmatrix} \in G'$ as in (A.3). By
\n(A.6), we see that $\mu(g') = \begin{pmatrix} \alpha & O \\ O & {}^t\alpha^{-1} \end{pmatrix}$. Now $g'^{-1} = \begin{pmatrix} {}^t\overline{\alpha} & -{}^t\beta \\ -\overline{t}\overline{\beta} & {}^t\alpha \end{pmatrix}$, so
\n
$$
\mu(g'^{-1}) = \begin{pmatrix} {}^t\overline{\alpha} & O \\ O & \overline{\alpha}^{-1} \end{pmatrix}.
$$
\n(A.7)

732 Andrew Knightly and Charles Li

A.4 Explicit formula for the matrix coefficient

Although it is possible to compute the matrix coefficient of π_{λ} explicitly for any λ (at least when $n = 2$), for simplicity, we have chosen here to treat just the case where dim $\Phi_{\lambda} = 1$. This discussion is valid for any $n \geq 2$.

For $k \in \mathbb{Z}$, let Φ_k be the character of $K = U(n)$ defined by

$$
\Phi_{\mathbf{k}}\begin{pmatrix} A & B \\ -B & A \end{pmatrix} = \det(A + Bi)^{\mathbf{k}}.
$$
\n(A.8)

Its holomorphic extension to $U_n(C) = GL_n(C)$ is given by the same formula. The (unique) weight of this character is

$$
\lambda_{k} = -k(\varepsilon_{1} + \dots + \varepsilon_{n}). \tag{A.9}
$$

This weight satisfies condition [\(A.4\)](#page-54-1) exactly when $k > n$.

PROPOSITION A.1. The character $\Phi_{\mathbf{k}} = \det^{\mathbf{k}} a$ rises as the minimal K-type of a holomorphic discrete series representation π_k^+ of $\text{Sp}_{2n}(\mathbf{R})$ if and only if $k > n$. If this holds and w_k is a unit vector of weight λ_k , then for any $g = \begin{pmatrix} A & B \\ C & D \end{pmatrix} \in$ $Sp_{2n}({\bf R})$ we have

$$
\left\langle \pi_{\mathbf{k}}^{+}(g)w_{\mathbf{k}},w_{\mathbf{k}}\right\rangle =\frac{2^{n\mathbf{k}}}{\det(A+D+i(-B+C))^{\mathbf{k}}}.
$$

Proof. By $(A.5)$ and $(A.7)$,

$$
\left\langle \pi_\mathbf{k}^+(g)w_\mathbf{k},w_\mathbf{k}\right\rangle=\Phi_\mathbf{k}(\mu(g^{-1}))^{-1}=\det({}^t\overline{\alpha})^{-\mathbf{k}}.
$$

By [\(A.3\)](#page-53-0), if $g = \begin{pmatrix} A & B \\ C & D \end{pmatrix}$,

$$
\alpha = \frac{1}{2}((A+D) + i(B-C)).
$$
 (A.10)

 \Box

The given formula now follows.

COROLLARY A.2. For $g = \begin{pmatrix} A & B \\ C & D \end{pmatrix} \in \text{Sp}_{2n}(\mathbf{R})$, define $f_{\mathbf{k}}^+(g) = \overline{\langle \pi_{\mathbf{k}}(g)w_{\mathbf{k}}, w_{\mathbf{k}} \rangle}$. Then

$$
f_{\mathbf{k}}^{+}(g) = \frac{2^{n\mathbf{k}}}{\det(A+D+i(B-C))^{\mathbf{k}}}.
$$
 (A.11)

COROLLARY A.3. For $k > n$ and $g = \begin{pmatrix} A & B \\ C & D \end{pmatrix} \in \text{Sp}_{2n}(\mathbf{R}), |f_k^+(g)|$ equals

$$
\frac{2^{nk}}{\det(2I_n + A^t A + B^t B + C^t C + D^t D + i(A^t C - C^t A + B^t D - D^t B))^{k/2}}.
$$

Proof. For $\alpha = A + D + i(B - C)$, $|\det \alpha|^2 = \det \alpha \det^{\dagger} \overline{\alpha} = \det \alpha^{\dagger} \overline{\alpha}$. Expand this using the relations given in (3.1) and (3.2) . The corollary then follows immediately. \Box

With the formula for the matrix coefficient in hand, we can compute its L^p norms. The formulas above and the calculations below closely parallel those for $GL_2(\mathbf{R})$ given in [\[KL1,](#page-68-4) §14].

PROPOSITION A.4. For any real number $\ell > 0$, the function $|f_{\mathbf{k}}^{+}|^{\ell}$ is integrable over $G = Sp_{2n}(\mathbf{R})$ if and only if $\ell \mathbf{k} > 2n$. If this condition holds, then with Haar measure normalized as in the proof below,

$$
\int_G |f_{\mathbf{k}}^+(g)|^\ell dg = \frac{2^{n(n+1)} \prod_{j=1}^{n-1} j!}{\prod_{1 \le i \le j \le n} (\ell \mathbf{k} - (i+j))}.
$$

Proof. The matrix coefficient $f_{\mathbf{k}}^{+}$ is bi-K-invariant, so it is convenient to use the Cartan decomposition $G = KA^+K$, where

$$
A^+ = \{a = \text{diag}(a_1, \dots, a_n, a_1^{-1}, \dots, a_n^{-1}) | 1 < a_1 < a_2 < \dots < a_n \}.
$$

We may view Δ^+ as the set of positive roots relative to the action of the diagonal subgroup on \mathfrak{g} , and $A^+ = \exp(\mathfrak{a}^+)$, where \mathfrak{a}^+ is the positive Weyl chamber. By a standard integration formula([\[vdB,](#page-70-0) Lemma 4.2]), when dg is suitably normalized we have

$$
\int_G |f_{\mathbf{k}}^+(g)|^{\ell} dg = \int_{K \times A^+ \times K} |f_{\mathbf{k}}^+(k_1ak_2)|^{\ell} \prod_{\alpha \in \Delta^+} (a^{\alpha} - a^{-\alpha}) dk_1 da dk_2,
$$

where the Haar measure of the compact group K is taken to be 1. Note that $(a^{\alpha} - a^{-\alpha}) > 0$ by the definition of A^{+} . We now change notation and write $a = diag(a_1, \ldots, a_n)$. Using Corollary [A.3,](#page-55-0) the above is

$$
= 2^{n\ell k} \int_{a \in GL_n(R) \text{ diagonal,}} \det(2I_n + a^2 + a^{-2})^{-\ell k/2} \prod_{\alpha \in \Delta^+} (a^{\alpha} - a^{-\alpha}) da
$$

$$
= 2^{n\ell k} \int_1^{\infty} \int_{a_1}^{\infty} \cdots \int_{a_{n-1}}^{\infty} \prod_{d=1}^n (2 + a_d^2 + a_d^{-2})^{-\ell k/2} (a_d^2 - a_d^{-2})
$$

$$
\times \prod_{1 \le i < j \le n} (a_i a_j - a_i^{-1} a_j^{-1}) (a_j a_i^{-1} - a_j^{-1} a_i) \frac{da_n}{a_n} \cdots \frac{da_1}{a_1}.
$$

To ease notation below, set $\kappa = \ell \mathbf{k}$. Then letting $u_j = 2 + a_j^2 + a_j^{-2}$, the above is

$$
=2^{n\kappa-n}\int_{4}^{\infty}\int_{u_1}^{\infty}\cdots\int_{u_{n-1}}^{\infty}\prod_{d=1}^{n}u_d^{-\kappa/2}\prod_{1\leq i (A.12)
$$

This can be evaluated using the Selberg integral, but we have chosen to give a self-contained treatment since it does not take much more space to do so. Observe that

$$
\prod_{i < j} (u_j - u_i) = \sum_{\sigma \in S_n} \operatorname{sgn}(\sigma) u_1^{\sigma(1)-1} \cdots u_n^{\sigma(n)-1}.
$$

Hence [\(A.12\)](#page-56-0) becomes

$$
\frac{2^{n\kappa}}{2^n} \sum_{\sigma \in S_n} \operatorname{sgn}(\sigma) \int_4^\infty \int_{u_1}^\infty \int_{u_2}^\infty \cdots \int_{u_{n-1}}^\infty u_1^{\sigma(1)-1-\kappa/2} \cdots u_n^{\sigma(n)-1-\kappa/2} du_n \cdots du_1
$$

$$
= \frac{2^{n\kappa}}{2^n} \sum_{\sigma \in S_n} \frac{\operatorname{sgn}(\sigma)(-1)^n 4^{\sigma(1)+\sigma(2)+\cdots+\sigma(n)-n\kappa/2}}{\prod_{j=1}^n (\sigma(n)+\sigma(n-1)+\cdots+\sigma(n-j+1)-j\kappa/2)}
$$

$$
= 2^{n^2} \sum_{\sigma \in S_n} \frac{\operatorname{sgn}(\sigma)}{(\kappa/2-\sigma(n)) (\kappa-\sigma(n)-\sigma(n-1))\cdots(n\kappa/2-\sigma(n)-\cdots-\sigma(1))},
$$

provided $\kappa > 2n$ (otherwise the integral diverges). Replace σ by $\sigma\tau$, where $\tau(i) = n - i + 1$ for all $1 \leq i \leq n$. The above is then

$$
=2^{n^2}\sum_{\sigma\in S_n}\frac{\operatorname{sgn}(\sigma\tau)}{(\kappa/2-\sigma(1))(\kappa-\sigma(1)-\sigma(2))\cdots(n\kappa/2-\sigma(1)-\cdots-\sigma(n))}.
$$

Applying Lemma [A.5](#page-57-0) below with $b_i = \kappa/2 - i$, the above is

$$
= 2^{n^2+n} \frac{\operatorname{sgn}(\tau) \prod_{i < j} (i-j)}{\prod_{i < j} (\kappa - (i+j))} = 2^{n(n+1)} \frac{\prod_{i < j} (j-i)}{\prod_{i < j} (\kappa - (i+j))},
$$

by the fact that $sgn(\tau) = (-1)^{\binom{n}{2}}$ (both are 1 iff $n \equiv 0, 1 \mod 4$). The proposition now follows.

LEMMA A.5. Let F be a field of characteristic 0, and let b_1, \ldots, b_n be indeterminate variables. Then in the field $F(b_1, \ldots, b_n)$ of rational functions,

$$
\sum_{\sigma \in S_n} \frac{\text{sgn}(\sigma)}{b_{\sigma(1)}(b_{\sigma(1)} + b_{\sigma(2)}) \cdots (b_{\sigma(1)} + \cdots + b_{\sigma(n)})} = 2^n \frac{\prod_{i < j} (b_i - b_i)}{\prod_{i \le j} (b_i + b_j)}.
$$
 (A.13)

Proof. Let $A(b_1, \ldots, b_n)$ denote the left-hand side of [\(A.13\)](#page-57-1). Then

$$
A(b_1,\ldots,b_n) = \sum_{\ell=1}^n (-1)^{n-\ell} \frac{A(b_1,\ldots,b_{\ell-1},b_{\ell+1},\ldots,b_n)}{b_1+\cdots+b_n}.
$$
 (A.14)

To see this, write $A(b_1, \ldots, b_n)$ as

$$
\frac{1}{b_1+\cdots+b_n}\sum_{\ell=1}^n\sum_{\sigma\in S_n,\atop \sigma(n)=\ell}\frac{\operatorname{sgn}(\sigma)}{b_{\sigma(1)}(b_{\sigma(1)}+b_{\sigma(2)})\cdots(b_{\sigma(1)}+\cdots+b_{\sigma(n-1)})}.
$$

Given $\sigma \in S_n$ with $\sigma(n) = \ell$, define $\sigma' \in S_{n-1} \subseteq S_n$ by

$$
\sigma'(i) = \begin{cases} \sigma(i) & \text{if } 1 \le \sigma(i) < \ell \\ \sigma(i) - 1 & \text{if } \ell + 1 \le \sigma(i) \le n. \end{cases}
$$

Then σ is the composition of σ' with $n - \ell$ transpositions, so sgn(σ) = $(-1)^{n-\ell}$ sgn(σ'). Since σ → σ' defines a bijection between the set of such σ and S_{n-1} , [\(A.14\)](#page-57-2) follows.

We may now prove [\(A.13\)](#page-57-1) by induction on n. The base case $n = 2$ is easy to check by hand. Applying the inductive hypothesis to [\(A.14\)](#page-57-2),

$$
A(b_1,\ldots,b_n) = \frac{2^{n-1}}{b_1 + \cdots + b_n} \sum_{\ell=1}^n (-1)^{n-\ell} \frac{\prod_{i < j; i, j \neq \ell} (b_j - b_i)}{\prod_{i \leq j; i, j \neq \ell} (b_i + b_j)}.
$$

Let

$$
B(b_1, ..., b_n) = \frac{b_1 + ... + b_n}{2^{n-1}} \left(\prod_{1 \le i \le j \le n} (b_i + b_j) \right) A(b_1, ..., b_n)
$$
 (A.15)

$$
= \sum_{\ell=1}^{n} (-1)^{n-\ell} \left(\prod_{i < j; i, j \neq \ell} (b_j - b_i) \right) \left(\prod_{i=1}^{n} (b_i + b_\ell) \right). \tag{A.16}
$$

This is a homogeneous polynomial of degree $\frac{(n-1)(n-2)}{2} + n = \frac{n(n-1)}{2} + 1$. Because

$$
A(b_1,\ldots,b_n)=\mathrm{sgn}(\sigma)A(b_{\sigma(1)},\ldots,b_{\sigma(n)})
$$

for all permutations $\sigma \in S_n$, B inherits this property from [\(A.15\)](#page-58-0). In particular,

$$
B(b_1,\ldots,b_n) = -B(b_{\sigma(1)},\ldots,b_{\sigma(n)})
$$
\n(A.17)

if $\sigma = (i \, j)$ is any 2-cycle. It follows that $B(b_1, \ldots, b_n) = 0$ if $b_i = b_j$ for any $i \neq j$. Hence $\prod_{i < j} (b_i - b_j)$ divides $B(b_1, \ldots, b_n)$, and

$$
\frac{B(b_1,\ldots,b_n)}{\prod_{i
$$

is a homogeneous symmetric polynomial of degree $\frac{n(n-1)}{2} + 1 - \frac{n(n-1)}{2} = 1$. Hence it has the form $c(b_1 + \cdots + b_n)$ for some constant c. The monomial $b_n^n \prod_{i=2}^{n-1} b_i^{i-1}$ appears in [\(A.16\)](#page-58-1) with coefficient 2 (take $\ell = n$), and in (b₁ + $\cdots + b_n$) $\prod_{i < j} (b_j - b_i)$ with coefficient 1. Therefore $c = 2$, and [\(A.13\)](#page-57-1) follows.

A.5 Integrability

An irreducible unitary representation π of $G = \mathrm{Sp}_{2n}(\mathbf{R})$ is said to be *integrable* if it has a nonzero matrix coefficient belonging to $L^1(G)$, or equivalently, if all of its K-finite matrix coefficients belong to $L^1(G)$. Applying Proposition [A.4](#page-56-1) with $\ell = 1$, we immediately obtain the following.

PROPOSITION A.6. The representation $\pi_{\mathbf{k}}^{+}$ is integrable if and only if $\mathbf{k} > 2n$.

More generally, let π_{λ} be the discrete series representation of G (holomorphic or not) with Harish-Chandra parameter λ . Then by a theorem due to Trombi, Varadarajan, Hecht and Schmid, π_{λ} is integrable if and only if

$$
|\langle \lambda, \beta \rangle| > \frac{1}{2} \sum_{\alpha \in \Delta^+} |\langle \alpha, \beta \rangle| \quad \text{for all} \quad \beta \in \Delta_{nc} \tag{A.18}
$$

([\[TV\]](#page-70-1), [\[HS\]](#page-68-5); see also Miličić [\[Mi\]](#page-69-0)). With notation as in $(A.2)$ and $(A.9)$, the Harish-Chandra parameter of $\pi_{\mathbf{k}}^{+}$ is

$$
\lambda = \lambda_{k} + \delta_{G} = (1 - k)\varepsilon_{1} + (2 - k)\varepsilon_{2} + \dots + (n - k)\varepsilon_{n}
$$
\n(A.19)

(see Remark (1) after Theorem 9.20 of [\[Kn\]](#page-68-3)). Note that for any j, ℓ ,

$$
|\langle \lambda_{\mathbf{k}} + \delta_{G}, \varepsilon_{j} + \varepsilon_{\ell} \rangle| = |(j + \ell) - 2\mathbf{k}| = 2\mathbf{k} - (j + \ell). \tag{A.20}
$$

Using this, one may easily verify that [\(A.18\)](#page-59-0) holds for λ exactly when $k > 2n$, for an alternative proof of Proposition [A.6.](#page-58-2)

A.6 Formal Degree

Recall that the formal degree of $\pi = \pi_{\lambda}$ is the constant $d_{\pi} > 0$ (depending only on the choice of Haar measure on $G = \mathrm{Sp}_{2n}(\mathbf{R}))$ satisfying

$$
\int_G |\langle \pi_\lambda(g)v, w\rangle|^2 dg = \frac{\|v\|^2 \|w\|^2}{d\lambda}
$$

for all $v, w \in V_{\pi}$. Applying Proposition [A.4](#page-56-1) with $\ell = 2$, we immediately find the following.

PROPOSITION A.7. The formal degree of π_k^+ is the following polynomial in k of degree $n + \binom{n}{2} = \frac{n(n+1)}{2}$ $\frac{i+1j}{2}$:

$$
d_{k} = a \prod_{1 \le i \le j \le n} (2k - (i + j)), \tag{A.21}
$$

where a is a nonzero constant depending on dg.

Remarks: (1) Harish-Chandra proved that there exists a choice of Haar measure for which

$$
d_\lambda = \prod_{\beta \in \Delta^+} \left| \frac{\langle \lambda + \delta_G, \beta \rangle}{\langle \delta_G, \beta \rangle} \right|
$$

for all λ ([\[HC\]](#page-67-4) §10). If λ is given by [\(A.19\)](#page-59-1), then by evaluating the above expression explicitly as in [\(A.20\)](#page-59-2), we obtain an alternative proof of [\(A.21\)](#page-59-3). (2) With measure normalized as in the proof of Proposition [A.4,](#page-56-1) $a =$ $(2^{n(n+1)} \prod_{j=1}^{n-1} j!)^{-1}.$

(3) If we adopt the classical normalization of measure, so that

$$
d_{\mathbf{k}}^{-1} = \int_G |f_{\mathbf{k}}^+(g)|^2 dg = \int_{\mathcal{H}_n} \left| f_{\mathbf{k}}^+\left(\left(\begin{smallmatrix} I_n & X \\ & I_n \end{smallmatrix} \right) \left(\begin{smallmatrix} Y^{1/2} \\ & Y^{-1/2} \end{smallmatrix} \right) \right) \right|^2 \frac{dXdY}{(\det Y)^{n+1}},
$$

then $a = 2^{-n(n+2)} \pi^{-n(n+1)/2}$ ([\[PSS,](#page-69-1) §A.1]).

A.7 EXTENSION OF $\pi_{\mathbf{k}}^{+}$ to GSp_{2n}

We can extend π_k^+ to a representation of $\mathrm{GSp}_{2n}(\mathbf{R})$ in the following way. First induce $\pi_{\mathbf{k}}^{+}$ to the group of symplectic similitudes with multiplier ± 1 , namely

$$
\mathrm{Sp}_{2n}^{\pm} = \left\langle \begin{pmatrix} I_n & \\ & -I_n \end{pmatrix} \right\rangle \mathrm{Sp}_{2n} .
$$

Let V^+ be the space of $\pi^+_{\mathbf{k}}$. Then the new space is

$$
V = \{ f : \operatorname{Sp}_{2n}^{\pm}(\mathbf{R}) \to V^+ | f(gx) = \pi_{\mathbf{k}}^+(g) f(x) \text{ for all } g \in \operatorname{Sp}_{2n}(\mathbf{R}) \},
$$

and $Sp_{2n}^{\pm}(\mathbf{R})$ acts by right translation. Note that any $f \in V$ is determined by $f(\begin{pmatrix} I_n & I_n \end{pmatrix})$ and $f(\begin{pmatrix} I_n & I_n \end{pmatrix})$. We identify V^+ with the subspace of V^{\pm} consisting of functions which vanish on $\binom{I_n}{I_n}$. Letting V^- denote the space of functions vanishing on the identity element, we have

$$
V = V^+ \oplus V^-.
$$
\n(A.22)

We make V into a Hilbert space by defining

$$
\langle f, h \rangle = \langle f(1), h(1) \rangle_{V^+} + \langle f(\sigma), h(\sigma) \rangle_{V^+}, \tag{A.23}
$$

where $\sigma =$ $\sqrt{I_n}$ $-I_n$ \setminus . Then [\(A.22\)](#page-60-0) is an orthogonal direct sum.

Denote this representation on V by π_k . One easily sees that each subspace in [\(A.22\)](#page-60-0) is stable under $\text{Sp}_{2n}(\mathbf{R})$, and $\pi_{\mathbf{k}}|_{\text{Sp}_{2n}(\mathbf{R})} = \pi_{\mathbf{k}}^+ \oplus \pi_{\mathbf{k}}^-$, where $\pi_{\mathbf{k}}^-$ is also irreducible and square integrable. Let

$$
Z^+ = \{ \begin{pmatrix} zI_n & \\ & zI_n \end{pmatrix} \mid z > 0 \}.
$$

Then

$$
GSp_{2n} = Z^+ \times Sp_{2n}^{\pm} .
$$

Extend $\pi_{\mathbf{k}}$ to a representation of $GSp_{2n}(\mathbf{R})$ by requiring Z^+ to act trivially. This is an irreducible square integrable representation, also denoted π_k . For any $z \in Z(\mathbf{R})$ and $f \in V$, we have

$$
\pi_{\mathbf{k}}(z)f(g) = f(\operatorname{sgn}(z)g) = \pi_{\mathbf{k}}^{+}(\operatorname{sgn}(z))f(g) = \operatorname{sgn}(z)^{n\mathbf{k}}f(g)
$$

by [\(A.8\)](#page-55-2). This shows that the central character of π_k is

$$
\chi_{\pi_{\mathbf{k}}}(z) = \text{sgn}(z)^{n\mathbf{k}}.\tag{A.24}
$$

As before, let $w_{\mathbf{k}} \in V^+$ denote a unit vector of weight $\lambda_{\mathbf{k}}$. Define $\phi_0 \in V$ by

$$
\phi_0(\begin{pmatrix} I_n & 0 \end{pmatrix}) = w_k
$$
 and $\phi_0(\begin{pmatrix} I_n & 0 \end{pmatrix}) = 0.$

This is a lowest weight vector, spanning the minimal K-type of π_k , which is the two-dimensional representation $\mathrm{Ind}_{K}^{K^{\pm}}(\Phi_{k}),$ where $K^{\pm} = K \cup K\begin{pmatrix} I_{n} & 0 \end{pmatrix}$.

PROPOSITION A.8. The representation π_k is irreducible, unitary, and squareintegrable when $k > n$. It is integrable exactly when $k > 2n$, and in this case, the formal degree of $\pi_{\mathbf{k}}$ coincides with that of $\pi_{\mathbf{k}}^{+}$ given in [\(A.21\)](#page-59-3).

Proof. Everything follows more or less immediately from the corresponding properties of $\pi^+_{\mathbf{k}}$. Indeed, define the matrix coefficient

$$
f_{\mathbf{k}}(g) = \overline{\langle \pi_{\mathbf{k}}(g)\phi_0, \phi_0 \rangle}.
$$
 (A.25)

Unraveling the definitions, we see that $f_{\mathbf{k}}(g) = f_{\mathbf{k}}^{+}(g)$ for $g \in \mathrm{Sp}_{2n}(\mathbf{R})$, and $f_{\mathbf{k}}(g) = 0$ if $r(g) < 0$. Using the fact that $\pi_{\mathbf{k}}$ is Z^+ -invariant, we have

$$
\int_{\mathrm{GSp}_{2n}(\mathbf{R})/Z} |f_{\mathbf{k}}(g)|^{\ell} dg = \int_{\mathrm{Sp}_{2n}^{\pm}} |f_{\mathbf{k}}(g)|^{\ell} dg = \int_{\mathrm{Sp}_{2n}} |f_{\mathbf{k}}(g)|^{\ell} dg = \int_{\mathrm{Sp}_{2n}} |f_{\mathbf{k}}^+(g)|^{\ell} dg.
$$

The assertions now follow from Proposition A.4.

The assertions now follow from Proposition [A.4.](#page-56-1)

THEOREM A.9. For
$$
g = \begin{pmatrix} A & B \\ C & D \end{pmatrix} \in \text{GSp}_{2n}(\mathbf{R}),
$$

$$
\langle \pi_{\mathbf{k}}(g)\phi_0, \phi_0 \rangle = \begin{cases} \frac{r(g)^{\frac{nk}{2}}2^{nk}}{\det(A+D+i(-B+C))^k} & \text{if } r(g) > 0 \\ 0 & \text{if } r(g) < 0. \end{cases}
$$

Proof. Let $r = r(g)$. Suppose $r < 0$. Then $\pi_k(g)\phi_0(x) \neq 0 \iff r(x) < 0$, i.e. $\pi_{\mathbf{k}}(g)\phi_0 \in V^-$, so it is orthogonal to $\phi_0 \in V^+$. Thus $\langle \pi_{\mathbf{k}}(g)\phi_0, \phi_0 \rangle = 0$ in this case.

Now suppose $r > 0$. Let

$$
h = r^{-1/2} g = \begin{pmatrix} r^{-1/2} A & r^{-1/2} B \\ r^{-1/2} C & r^{-1/2} D \end{pmatrix}.
$$

It is easy to see that $r(h) = 1$, i.e. that $h \in \text{Sp}_{2n}(\mathbf{R})$. By definition of π_k , $\pi_{\mathbf{k}}(g)\phi_0 = \pi_{\mathbf{k}}(h)\phi_0$. Therefore by [\(A.23\)](#page-60-1),

$$
\langle \pi_{\mathbf{k}}(g)\phi_0, \phi_0 \rangle = \langle \pi_{\mathbf{k}}(h)\phi_0, \phi_0 \rangle = \langle \pi_{\mathbf{k}}^+(h)w_{\mathbf{k}}, w_{\mathbf{k}} \rangle_{V^+}
$$

$$
= \frac{2^{nk}}{r^{-nk/2}\det(A+D+i(-B+C))^k}
$$

by Proposition [A.1.](#page-55-3)

The following corollaries are easily proven as in §[A.4.](#page-55-4) COROLLARY A.10. For $f_{\mathbf{k}}(g) = \overline{\langle \pi_{\mathbf{k}}(g)\phi_0, \phi_0 \rangle}$, if $r(g) > 0$, then

$$
f_{k}(g) = \frac{r(g)^{\frac{nk}{2}} 2^{nk}}{\det(A+D+i(B-C))^{k}}.
$$

If $r(q) < 0$, then $f_k(q) = 0$.

COROLLARY A.11. If $g = \begin{pmatrix} A & B \\ C & D \end{pmatrix} \in \text{GSp}_{2n}(\mathbf{R})$ with $r(g) > 0$, then $|f_{k}(g)|$ equals

$$
\frac{r(g)^{\frac{n\mathbf{k}}{2}}2^{n\mathbf{k}}}{\det(2r(g)I_n + A^tA + B^tB + C^tC + D^tD + i(A^tC - C^tA + B^tD - D^tB))^{k/2}}.
$$

In the special case $n = 2$, we can make the above more explicit and provide a convenient upper bound for the matrix coefficient.

PROPOSITION A.12. Suppose $n = 2$, and $r(g) > 0$. Then

$$
|f_{\mathbf{k}}(g)|=\frac{r(g)^{\mathbf{k}}\,4^{\mathbf{k}}}{(4r(g)^2+2r(g)\sum_{i,j}g_{ij}^2+\sum_{i=3}^8X_i^2)^{\mathbf{k}/2}},
$$

where g_{ij} are the entries of g, and the X_i are the bilinear forms in these entries defined in the proof below. Consequently,

$$
|f_{k}(g)| \le \frac{(8r(g))^{k/2}}{(2r(g) + \sum_{i,j} g_{ij}^{2})^{k/2}}.
$$
\n(A.26)

Proof. Write $g = \begin{pmatrix} A & B \\ C & D \end{pmatrix}$. Let A_{ij} denote the (i, j) -th entry of A, and likewise for the entries of B, C, D . Then

$$
det(2r(g)I_2 + A^t A + B^t B + C^t C + D^t D + i(A^t C - C^t A + B^t D - D^t B))
$$

= $(2r(g) + A_{11}^2 + A_{12}^2 + B_{11}^2 + B_{12}^2 + B_{11}^2 + C_{12}^2 + D_{11}^2 + D_{12}^2)$ (A.27)
 $\times (2r(g) + A_{21}^2 + A_{22}^2 + B_{21}^2 + B_{22}^2 + C_{21}^2 + C_{22}^2 + D_{21}^2 + D_{22}^2) - X_1^2 - X_2^2,$

where

$$
X_1 = A_{11}A_{21} + A_{12}A_{22} + B_{11}B_{21} + B_{12}B_{22} + C_{11}C_{21} + C_{12}C_{22} + D_{11}D_{21} + D_{12}D_{22},
$$

\n
$$
X_2 = A_{11}C_{21} + A_{12}C_{22} + B_{11}D_{21} + B_{12}D_{22} - C_{11}A_{21} - C_{12}A_{22} - D_{11}B_{21} - D_{12}B_{22}.
$$

\nBy Degen's eight-square identity,

$$
(A_{11}^2 + A_{12}^2 + B_{11}^2 + B_{12}^2 + C_{11}^2 + C_{12}^2 + D_{11}^2 + D_{12}^2)
$$

\$\times (A_{21}^2 + A_{22}^2 + B_{21}^2 + B_{22}^2 + C_{21}^2 + C_{22}^2 + D_{21}^2 + D_{22}^2) = \sum_{i=1}^8 X_i^2\$,

for X_1, X_2 as above, and

$$
X_3 = A_{11}A_{22} - A_{12}A_{21} + B_{11}D_{22} - B_{12}D_{21} - C_{11}C_{22} + C_{12}C_{21} - D_{12}B_{21} + D_{11}B_{22},
$$

\n
$$
X_4 = A_{11}C_{22} - A_{12}C_{21} + B_{11}B_{22} - B_{12}B_{21} - C_{12}A_{21} + C_{11}A_{22} - D_{11}D_{22} + D_{12}D_{21},
$$

\n
$$
X_5 = A_{11}B_{21} - A_{12}D_{22} - B_{11}A_{21} + B_{12}C_{22} - C_{11}D_{21} - C_{12}B_{22} + D_{12}A_{22} + D_{11}C_{21},
$$

\n
$$
X_6 = A_{11}D_{21} - A_{12}B_{22} - B_{11}C_{21} + B_{12}A_{22} + C_{11}B_{21} + C_{12}D_{22} - D_{11}A_{21} - D_{12}C_{22},
$$

 $X_7 = A_{11}D_{22} + A_{12}B_{21} - B_{11}A_{22} - B_{12}C_{21} + C_{11}B_{22} - C_{12}D_{21} + D_{11}C_{22} - D_{12}A_{21},$ $X_8 = A_{11}B_{22} + A_{12}D_{21} - B_{11}C_{22} - B_{12}A_{21} - C_{11}D_{22} + C_{12}B_{21} - D_{11}A_{22} + C_{21}D_{12}.$ Therefore [\(A.27\)](#page-62-0) equals

$$
4r(g)^{2} + 2r(g)\sum_{i,j} g_{ij}^{2} + \sum_{i=3}^{8} X_{i}^{2},
$$

and the proposition follows from Corollary [A.11.](#page-62-1)

B Off-diagonal terms

In this appendix, we give a very rough estimate for the off-diagonal terms in the relative trace formula to obtain the following quantitative version of Theorem [6.3](#page-24-0) for a fixed level N, valid when $n = 2$.

THEOREM B.1. Suppose $n = 2$, $k \ge 17$, S is a finite set of primes, $r = \prod_{p \in S} p^{r_p}$ is the similitude attached to the local test functions f_p as in [\(5.1\)](#page-15-0), and N is a fixed level prime to S . Then with notation as in Theorem [6.3,](#page-24-0)

$$
\frac{1}{\psi(N)} \sum_{\pi \in \Pi_{\mathbf{k}}(N)} \sum_{\varphi \in E_{\mathbf{k}}(\pi, N)} \frac{c_{\sigma_1}(\varphi) c_{\sigma_2}(\varphi)}{\|\varphi\|^2} \prod_{p \in \mathbf{S}} (\mathcal{S}f_p)(t_{\pi_p})
$$
\n
$$
= \sum_A \int_{S_2(\mathbf{A})} f_1(\begin{pmatrix} I_2 & S \\ O & I_2 \end{pmatrix} \begin{pmatrix} A & O \\ O & r^t A^{-1} \end{pmatrix}) \theta(\text{tr } \sigma_1 S) dS + O\left(\frac{\mathbf{k}^{21/2} (8r)^{\mathbf{k}/2}}{N^{\mathbf{k}-12}}\right),
$$

for an absolute implied constant.

Remarks: (1) The sum over A is dependent on k , but not on N .

(2) With extra work, one can increase the power of N in the error term, and decrease the lower bound on k.

(3) As noted, this is a very rough estimate, which is meaningful only when k is fixed. Indeed in the k-aspect, the "error term" is actually larger than the "main term" when k is large, as the latter contains the factor $c_{n\mathbf{k}\sigma_1}$ of Theorem [1.1](#page-3-0) with very rapid decay in k.

Before specializing to the case $n = 2$, we give the general form of the Fourier trace formula for a fixed level N and any n. With notation as in §[5-](#page-14-0)[6,](#page-20-0) define

$$
I_N = \frac{1}{\psi(N)} \iint\limits_{(N(\mathbf{Q}) \backslash N(\mathbf{A}))^2} K_{f_N}(n_1, n_2) \overline{\theta_1(n_1)} \theta_2(n_2) dn_1 dn_2.
$$

By using the spectral and geometric forms of the kernel, we have

$$
\frac{1}{\psi(N)} \sum_{\pi \in \Pi_{\mathbf{k}}(N)} \sum_{\varphi \in E_{\mathbf{k}}(\pi, N)} \frac{c_{\sigma_1}(\varphi) c_{\sigma_2}(\varphi)}{\|\varphi\|^2} \prod_{p \in \mathbf{S}} (\mathcal{S}f_p)(t_{\pi_p}) = M(f) + E(f), \quad (B.1)
$$

Documenta Mathematica 24 (2019) 677–747

 \Box

where

$$
M(f) = \sum_{A} \int_{S_n(A)} f_1(\begin{pmatrix} I_n & S \\ O & I_n \end{pmatrix} \begin{pmatrix} A & O \\ O & r^t A^{-1} \end{pmatrix}) \theta(\text{tr } \sigma_1 S) dS
$$

is the "diagonal" term appearing in Theorem [6.3,](#page-24-0) and

$$
E(f) = \frac{1}{\psi(N)} \iint_{(N(\mathbf{Q}) \backslash N(\mathbf{A}))^2} \sum_{\gamma = \begin{pmatrix} A & B \\ C & D \end{pmatrix} \in \mathbb{G}(\mathbf{Q}), C \neq O} f(n_1^{-1} \gamma n_2) \overline{\theta_1(n_1)} \theta_2(n_2) dn_1 dn_2
$$

is the off-diagonal contribution. (One checks readily that the integrand is indeed invariant under $N(Q) \times N(Q)$.)

It is possible to express $E(f)$ as a sum of explicit orbital integrals (cf. [\[KL2,](#page-68-6) §2]), with the orbits determined via the Bruhat decomposition of $\mathbb{G}(\mathbf{Q})$. For our more modest goal of proving Theorem [B.1,](#page-63-0) it suffices to show that

$$
E(f) \ll \frac{\mathbf{k}^{21/2} (8r)^{\mathbf{k}/2}}{N^{\mathbf{k}-12}} \tag{B.2}
$$

when $n = 2$.

LEMMA B.2. For $k \ge 2$, $\Delta > 0$, and $a \in \mathbb{R}$,

$$
\sum_{n=-\infty}^{\infty} \frac{1}{((n+a)^2 + \Delta^2)^{k/2}} \le \frac{k+\Delta}{\Delta} \int_{\mathbf{R}} \frac{dx}{(x^2 + \Delta^2)^{k/2}}.
$$

Proof. Let

 ϵ

$$
f(x) = \frac{1}{((x+a)^2 + \Delta^2)^{k/2}}.
$$

Then $f'(x) = -\frac{k(x+a)}{(x+a)^2 + \Delta^2}$ $\frac{f(x+a)}{((x+a)^2+\Delta^2)^{k/2+1}}$, so that

$$
|f'(x)| \le k \frac{((x+a)^2 + \Delta^2)^{1/2}}{((x+a)^2 + \Delta^2)^{k/2+1}} = \frac{k}{((x+a)^2 + \Delta^2)^{1/2}} f(x) \le \frac{k}{\Delta} f(x).
$$

By Euler's summation formula([\[MV,](#page-69-2) p. 495]),

$$
\sum_{\substack{
$$

where $\{x\} = x - [x]$ is the fractional part of x. Hence

$$
\sum_{n=-\infty}^{\infty} f(n) \le \int_{\mathbf{R}} f(x)dx + \int_{\mathbf{R}} |f'(x)|dx \le (\frac{\mathbf{k}}{\Delta} + 1) \int_{\mathbf{R}} f(x)dx
$$

$$
= \frac{\mathbf{k} + \Delta}{\Delta} \int_{\mathbf{R}} \frac{dx}{((x+a)^2 + \Delta^2)^{k/2}} = \frac{\mathbf{k} + \Delta}{\Delta} \int_{\mathbf{R}} \frac{dx}{(x^2 + \Delta^2)^{k/2}}.
$$

Proof of Theorem [B.1.](#page-63-0) As before, for any set R, we let $S_2(R)$ denote the 2×2 symmetric matrices over R. For such a matrix S, let $n_S = \begin{pmatrix} I_2 & S \\ O & I_2 \end{pmatrix}$. Because $[0, 1) \times \hat{\mathbf{Z}}$ is a fundamental domain for $\mathbf{Q} \setminus \mathbf{A}$,

$$
E(f) = \frac{1}{\psi(N)} \iint_{S_2([0,1)\times \widehat{\mathbf{Z}})^2_{\gamma = \left(\begin{smallmatrix} A & B \\ C & D \end{smallmatrix}\right) \in \mathbb{G}(\mathbf{Q}), C \neq O} f(n_S^{-1} \gamma n_{S'}) \theta(\operatorname{tr} \sigma_2 S' - \operatorname{tr} \sigma_1 S) \, dS dS'.
$$

Assuming $n_{S_{\text{fin}}}^{-1} \gamma n_{S'_{\text{fin}}} \in \text{supp } f_{\text{fin}}$, by Proposition [5.1](#page-15-1) we may take

$$
r(n_S^{-1}\gamma n_{S'})=r(\gamma)=\pm r,\quad (n_S^{-1}\gamma n_{S'})_{\text{fin}}\in M_4(\widehat{\mathbf{Z}}),\quad C\in M_2(N\mathbf{Z}).
$$

In fact, because f_{∞} is supported on matrices with positive similitude, we can take $r(\gamma) = r$. By the fact that S_{fin} , $S'_{fin} \in S_2(\mathbf{Z})$, it also follows that $A, B, D \in \mathbb{Z}$ $M_2(\mathbf{Z})$.

For a square matrix q with real entries, define

$$
\mathcal{Q}(g) = \sum_{i,j} g_{ij}^2
$$

where g_{ij} is the (i, j) -th entry of g. Then since $|f(n_1^{-1}\gamma n_2)| \le$ $\psi(N)|f_{\infty}(n_{1,\infty}^{-1}\gamma n_{2,\infty})|$, it follows by [\(A.26\)](#page-62-2) that

$$
|E(f)| \leq \frac{1}{2} d_{\mathbf{k}}(8r)^{\mathbf{k}/2} \iint_{S_2([0,1))^2} \sum_{\substack{A,B,D \in \mathbf{M}_2(\mathbf{Z}),\\O \neq C \in M_2(N\mathbf{Z})}} (2r + \mathcal{Q}(n_S^{-1} \begin{pmatrix} A & B \\ C & D \end{pmatrix} n_{S'}))^{-\mathbf{k}/2} dS dS'.
$$
\n(B.3)

(The factor $1/2$ accounts for the fact that we quotient by the center of $G(\mathbf{Q})$.) We first consider the sum over B :

$$
\sum_{B \in M_2(\mathbf{Z})} \left(2r + \mathcal{Q}(n_S^{-1} \begin{pmatrix} A & B \\ C & D \end{pmatrix} n_{S'})\right)^{-k/2}
$$

=
$$
\sum_{B \in M_2(\mathbf{Z})} \left(2r + \mathcal{Q}(\begin{pmatrix} A - SC & B - SD + AS' - SCS' \\ C & CS' + D \end{pmatrix})\right)^{-k/2}.
$$

By four applications of Lemma [B.2,](#page-64-0) the above is

$$
\ll k^4 \int_{M_2(\mathbf{R})} \left(2r + \mathcal{Q}(\begin{pmatrix} A - SC & Y \\ C & CS' + D \end{pmatrix})\right)^{-k/2} dY.
$$

Summing this over A, D , we find in the same way that

$$
\sum_{A,B,D \in M_2(\mathbf{Z})} \left(2r + \mathcal{Q}(n_S^{-1} \begin{pmatrix} A & B \\ C & D \end{pmatrix} n_{S'})\right)^{-k/2}
$$

$$
\ll \kappa^{12} \int_{M_2(\mathbf{R})^3} \left(2r + \mathcal{Q}(\begin{pmatrix} X & Y \\ C & Z \end{pmatrix})\right)^{-k/2} dX dY dZ.
$$

We now have an expression that is independent of S, S' , so using $(A.21), (B.3)$ $(A.21), (B.3)$ gives

$$
|E(f)|\ll \mathbf{k}^{15}(8r)^{\mathbf{k}/2}\sum_{O\neq C\in M_2(N{\bf Z})}\int_{M_2({\bf R})^3}\left(2r+\mathcal{Q}(\begin{pmatrix}X&Y\\C&Z\end{pmatrix})\right)^{-\mathbf{k}/2}dXdYdZ.
$$

Because $C \neq O$, the above integral is bounded above by

$$
\int_{M_2(\mathbf{R})^3} \mathcal{Q}(\begin{pmatrix} X & Y \ C & Z \end{pmatrix})^{-k/2} dX dY dZ
$$

=
$$
\int_{M_2(\mathbf{R})^3} (\mathcal{Q}(X) + \mathcal{Q}(Y) + \mathcal{Q}(Z) + \mathcal{Q}(C))^{-k/2} dX dY dZ.
$$

Replacing X, Y, Z with $Q(C)^{1/2}X$, $Q(C)^{1/2}Y$, $Q(C)^{1/2}Z$ respectively, the above is

$$
= \int_{\mathbf{R}^{12}} \frac{\mathcal{Q}(C)^{-k/2+6} dX_{11} \cdots dX_{22} dY_{11} \cdots dY_{22} dZ_{11} \cdots dZ_{22}}{(1+X_{11}^2+\cdots+X_{22}^2+Y_{11}^2+\cdots+Y_{22}^2+Z_{11}^2+\cdots+Z_{22}^2)^{k/2}},
$$

which converges if $k \geq 13$. Indeed, for such k, we find using spherical coordinates that the above is

$$
\ll \mathcal{Q}(C)^{-\mathbf{k}/2+6} \int_0^\infty \frac{\rho^{11}}{(1+\rho^2)^{\mathbf{k}/2}} d\rho \ll \mathbf{k}^{-6} (C_{11}^2 + C_{12}^2 + C_{21}^2 + C_{22}^2)^{-\mathbf{k}/2+6}
$$

since the integral over ρ is equal to $\frac{1}{2}B(6, \frac{k}{2} - 6)$ for the Beta function $B(x, y)$, and by Stirling's formula, $B(6, \frac{k}{2} - \bar{6}) \sim \Gamma(6)(\frac{k}{2} - 6)^{-6}$ as $k \to \infty$.

Finally, we need to sum over all $O \neq C \in M_2(N{\bf Z})$. Write $C_{ij} = NC'_{ij}$. We may assume $C_{11} \neq 0$. (The other cases can be handled by the exactly the same method.) Thus

$$
|E(f)| \ll \frac{\mathbf{k}^{9}(8r)^{k/2}}{N^{k-12}} \sum_{C'_{ij} \in \mathbf{Z}, C'_{11} \neq 0} ((C'_{11})^{2} + (C'_{12})^{2} + (C'_{21})^{2} + (C'_{22})^{2})^{-k/2+6}.
$$

$$
\ll \frac{\mathbf{k}^{12}(8r)^{k/2}}{N^{k-12}} \sum_{C'_{11} \neq 0} \int_{\mathbf{R}^{3}} ((C'_{11})^{2} + x^{2} + y^{2} + z^{2})^{-k/2+6} dx dy dz \text{ (by Lemma B.2)}
$$

$$
= \frac{\mathbf{k}^{12}(8r)^{k/2}}{N^{k-12}} \sum_{c \neq 0} c^{-k+15} \int_{\mathbf{R}^{3}} (1 + x^{2} + y^{2} + z^{2})^{-k/2+6} dx dy dz
$$

$$
\ll \frac{\mathbf{k}^{12}(8r)^{k/2}}{N^{k-12}} \sum_{c \neq 0} \frac{1}{c^{k-15}} \int_{0}^{\infty} \frac{\rho^{2}}{(1 + \rho^{2})^{k/2-6}} d\rho.
$$

This converges as long as $k \geq 17$. The integral is equal to $\frac{1}{2}B(\frac{3}{2}, \frac{k-15}{2}) \sim$ $\frac{1}{2}\Gamma(\frac{3}{2})(\frac{k-15}{2})^{-3/2} \ll k^{-3/2}$. Hence for $k \ge 17$, the above is $\ll \frac{k^{21/2}(8r)^{k/2}}{N^{k-12}}$ for an absolute implied constant. This proves [\(B.2\)](#page-64-1), and Theorem [B.1](#page-63-0) follows.

REFERENCES

- [AS] M. Asgari and R. Schmidt, Siegel modular forms and representations, Manuscripta Math. 104 (2001), no. 2, 173–200.
- [BBR] V. Blomer, J. Buttcane, and N. Raulf, A Sato-Tate law for GL(3), Comm. Math. Helv. 89 (2014), no. 4, 895-919.
- [Bi] B. J. Birch, *How the number of points of an elliptic curve over a fixed* prime field varies, J. London Math. Soc., 43 (1968), 57-60.
- [Bl] \quad V. Blomer, Spectral summation formula for $GSp(4)$ and moments of spinor L-functions, J. Eur. Math. Soc., to appear. [\(arXiv:1602.00780\)](https://arxiv.org/abs/1602.00780).
- [Br] R. Bruggeman, Fourier coefficients of cusp forms, Invent. Math. 45 (1978), no. 1, 1-18.
- [BrM] R. Bruggeman and R. Miatello, Eigenvalues of Hecke operators on Hilbert modular groups, Asian J. Math. 17 (2013), no. 4, 729-757.
- [BW] A. Borel and N. Wallach, *Continuous cohomology, discrete subgroups*, and representations of reductive groups, Second edition. Mathematical Surveys and Monographs, 67. Amer. Math. Soc., Providence, RI, 2000.
- [Bu] D. Bump, Lie Groups, Springer-Verlag, New York, 2004.
- [Ca] P. Cartier, Representations of p-adic groups: a survey, pp. 111-155, Proc. Sympos. Pure Math. XXXIII Part 1, Amer. Math. Soc., Providence, 1979.
- [CKM] M. Chida, H. Katsurada, and K. Matsumoto, On Fourier coefficients of Siegel modular forms of degree two with respect to congruence subgroups, Abh. Math. Semin. Univ. Hambg. 84 (2014), no. 1, 31-47.
- [D] M. Dickson, Local spectral equidistribution for degree 2 Siegel modular forms in level and weight aspects, Int. J. Number Theory 11 (2015), no. 2, 341-396.
- [FH] W. Fulton and J. Harris, Representation Theory, Springer-Verlag, New York, 1991.
- [Gr] B. Gross, On the Satake isomorphism. In: Galois representations in arithmetic algebraic geometry (Durham, 1996), 223-237, London Math. Soc. Lecture Note Ser., 254, Cambridge Univ. Press, Cambridge, 1998.
- [H] C. Herz, Bessel functions of matrix argument, Ann. of Math. (2) 61, (1955). 474-523.
- [HC] Harish-Chandra, Representations of semisimple Lie groups VI. Integrable and square-integrable representations, Amer. J. Math. 78 (1956), 564-628.

- [HKP] T. Haines, R. Kottwitz, and A. Prasad, Iwahori-Hecke algebras, J. Ramanujan Math. Soc., 25 (2010), no. 2, 113-145.
- [HS] H. Hecht and W. Schmid, On integrable representations of a semisimple Lie group, Math. Ann. 220 (1976), no. 2, 147-149.
- [Hua] Hua L. K., Introduction to number theory, Springer-Verlag, Berlin-New York, 1982.
- [I] A. E. Ingham, An integral which occurs in statistics, Proc. Cambridge Philos. Soc. 29, (1933), 1-26.
- [JS] H. Jacquet and J. Shalika, On Euler products and the classification of automorphic representations I & II, Amer. J. Math. 103 (1981), 499-588 & 777-815.
- [Ki] Y. Kitaoka, Fourier coefficients of Siegel cusp forms of degree two, Nagoya Math. J. 93 (1984), 149-171.
- [Kn] A. Knapp, Representation theory of semisimple groups, an overview based on examples, Princeton University Press, Princeton, NJ, 1986.
- [Kl] Klingen, H Introductory lectures on Siegel modular forms. Cambridge Studies in Advanced Mathematics, 20. Cambridge University Press, Cambridge, 1990.
- [KL1] A. Knightly and C. Li, Traces of Hecke operators, Mathematical Surveys and Monographs, 133. Amer. Math. Soc., 2006.
- [KL2] A. Knightly and C. Li, Petersson's trace formula and the Hecke eigenvalues of Hilbert modular forms, in "Modular forms on Schiermonnikoog," Cambridge University Press, 2008.
- [KL3] A. Knightly and C. Li, Kuznetsov's trace formula and the Hecke eigenvalues of Maass forms, Mem. Amer. Math. Soc. 224 (2013), no. 1055.
- [KR] A. Knightly and C. Reno, Weighted distribution of low-lying zeros of GL(2) L-functions, Canad. J. Math., to appear.
- [KST1] E. Kowalski, A. Saha, and J. Tsimerman, A note on Fourier coefficients of Poincaré series, Mathematika 57 (2011), no. 1, 31-40.
- [KST2] A. Saha, and J. Tsimerman, Local spectral equidistribution for Siegel modular forms and applications, Compos. Math. 148 (2012), no. 2, 335-384.
- [KWY] H. Kim, S. Wakatsuki, and T. Yamauchi, An equidistribution theorem for holomorphic Siegel modular forms for $GSp₄$ and its applications, preprint [\(arXiv:1604.02036\)](https://arxiv.org/abs/1604.02036), 2016.

746 Andrew Knightly and Charles Li

- [Li] C. Li, On the distribution of Satake parameters of GL_2 holomorphic cuspidal representations, Israel J. Math. 169 (2009), 341-373.
- [M] H. Maass, Über die Darstellung der Modulformen n-ten Grades durch Poincarésche Reihen, Math. Ann. 123, (1951), 125-151.
- [Mi] D. Miličić, Asymptotic behaviour of matrix coefficients of the discrete series, Duke Math. J. 44 (1977), no. 1, 59-88.
- [MT] J. Matz and N. Templier, Sato-Tate equidistribution for families of Hecke-Maass forms on $SL(n, \mathbf{R})/SO(n)$, preprint [\(arXiv:1505.07285\)](https://arxiv.org/abs/1505.07285), 2016.
- [MV] H. Montgomery and R. Vaughan, Multiplicative number theory I. Classical theory, Cambridge Studies in Advanced Mathematics, 97, Cambridge University Press, Cambridge, 2007.
- [Ne] M. Newman, Integral matrices, Pure and Applied Mathematics, Vol. 45. Academic Press, New York-London, 1972.
- [PSS] A. Pitale, A. Saha, and R. Schmidt, On the standard L-function for $GSp_{2n} \times GL_1$ and algebraicity of symmetric fourth L-values for $GL(2)$, preprint [\(arXiv:1803.06227\)](https://arxiv.org/abs/1803.06227), 2018.
- $[Ro]$ J. Rogawski, Representations of $GL(n)$ and division algebras over a p-adic field, Duke Math. J. 50 (1983), no. 1, 161-196.
- [Ru] W. Rudin, Functional Analysis, Second Edition, Int. Ser. Pure and App. Math., McGraw-Hill, Inc., New York, 1991.
- [Sa] A. Saha, On ratios of Petersson norms for Yoshida lifts, Forum Math. 27 (2015), no. 4, 2361-2412.
- [Sha] F. Shahidi, Eisenstein series and automorphic L-functions, Colloquium Publications, 58. American Mathematical Society, Providence, RI, 2010.
- [Shin] S. W. Shin, Automorphic Plancherel Density Theorem, Israel J. Math. 192 (2012), no. 1, 83-120.
- [Si1] C. Siegel, Uber die analytische Theorie der quadratischen Formen, Ann. of Math. (2) 36 (1935), no. 3, 527-606.
- [Si2] C. Siegel, Symplectic geometry, Amer. J. Math. 65, (1943). 1-86.
- [SST] P. Sarnak, S. W. Shin, and N. Templier, Families of L-functions and their symmetry, in: Families of automorphic forms and the trace formula, 531-578, Simons Symp., Springer, [Cham], 2016.

- [ST] S. W. Shin and N. Templier, Sato-Tate theorem for families and lowlying zeros of automorphic L-functions, Invent. Math. 203 (2016), no. 1, 1-177.
- [TV] P. Trombi and V. S. Varadarajan, Asymptotic behaviour of eigen functions on a semisimple Lie group: the discrete spectrum, Acta. Math. 129 (1972), no. 3-4, 237-280.
- [vdB] E. P. van den Ban, Induced representations and the Langlands classification, in Representation Theory and Automorphic Forms, pp. 123 - 155, T. N. Bailey, A. W. Knapp, eds, Proceedings of Symposia in Pure Mathematics, Vol. 61, Amer. Math. Soc. 1997.
- [W] F. Waibel, Moments of spinor L-functions and symplectic Kloosterman sums, preprint [\(arXiv:1805.11502\)](https://arxiv.org/abs/1805.11502), 2018.
- [Z] F. Zhou, Weighted Sato-Tate Vertical Distribution of the Satake Parameter of Maass Forms on $PGL(N)$, Ramanujan J. 35 (2014), no. 3, 405-425.

Andrew Knightly Department of Mathematics & Statistics University of Maine Neville Hall Orono, ME 04469-5752 USA knightly@math.umaine.edu

Charles Li Department of Mathematics The Chinese University of Hong Kong Shatin Hong Kong charles@charlesli.org

Documenta Mathematica 24 (2019)

748