ASYMPTOTICS OF HARISH-CHANDRA EXPANSIONS, BOUNDED HYPERGEOMETRIC FUNCTIONS ASSOCIATED WITH ROOT SYSTEMS, AND APPLICATIONS

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Abstract. A series expansion for Heckman-Opdam hypergeometric functions $\varphi_\lambda$ is obtained for all $\lambda \in a^*_\mathbb{C}$. As a consequence, estimates for $\varphi_\lambda$ away from the walls of a Weyl chamber are established. We also characterize the bounded hypergeometric functions and thus prove an analogue of the celebrated theorem of Helgason and Johnson on the bounded spherical functions on a Riemannian symmetric space of the noncompact type. The $L^p$-theory for the hypergeometric Fourier transform is developed for $0 < p < 2$. In particular, an inversion formula is proved when $1 \leq p < 2$.

Introduction

A natural extension of Harish-Chandra’s theory of spherical functions on Riemannian symmetric spaces of the noncompact type was introduced by Heckman and Opdam in the late eighties ([14], [12], [21]). In this theory, the symmetric space is replaced by a triple $(a, \Sigma, m)$ consisting of a finite dimensional real Euclidean vector space $a$, a root system $\Sigma$ in the dual $a^*$ of $a$, and a positive multiplicity function $m$ on $\Sigma$. A commuting family $D = D(a, \Sigma, m)$ of differential operators on $a$ is associated with this triple. The hypergeometric functions of Heckman and Opdam are joint eigenfunctions of $D$. For certain values of the multiplicity function, the triple $(a, \Sigma, m)$ indeed arises from a Riemannian symmetric space of the noncompact type $G/K$. In this case, $D$ coincides with the algebra of radial components of the $G$-invariant differential operators on $G/K$, and Heckman-Opdam’s hypergeometric functions are the restrictions to $a$ of Harish-Chandra’s elementary spherical functions on $G/K$. Heckman-Opdam’s theory of hypergeometric functions associated with root systems underwent an important development with the discovery of Cherednik operators (see [3], [22], [23] and references therein). The Cherednik operators (also called Dunkl-Cherednik operators or trigonometric Dunkl operators, as they are the curved analogue of the Dunkl operators on $\mathbb{R}^n$) are a commuting family of first order differential-reflection operators. They allow to construct algebraically all elements of $D$.

Let $W$ denote the Weyl group of $\Sigma$, and let $C^\infty_c(a)^W$ be the space of compactly-supported $W$-invariant smooth functions on $a$. The spectral decomposition of $D$ on $C^\infty_c(a)^W$ is obtained by means of the hypergeometric Fourier transform. Let $a^*_\mathbb{C}$ be the complexified dual of $a$, and

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let $\varphi_{\lambda}$ denote the Heckman-Opdam’s hypergeometric function of spectral parameter $\lambda$. The hypergeometric Fourier transform $\hat{f}$ (or $Ff$) of a sufficiently regular $W$-invariant function $f$ on $a$ is defined by integration against the $\varphi_{\lambda}$’s:

$$ (Ff)(\lambda) = \hat{f}(\lambda) = \int_a f(x)\varphi_{\lambda}(x)\,d\mu(x). $$

Here $d\mu(x)$ is a suitable measure on $a$ attached to the triple $(a, \Sigma, m)$ (see subsection 1.5). In [22], Opdam established the basic results in the $L^2$-harmonic analysis of the hypergeometric Fourier transform: the Paley-Wiener theorem, characterizing the image under $F$ of $C_c^\infty(a)^W$, the inversion formula on $C_c^\infty(a)^W$, and the Plancherel theorem. The $L^2$-Schwartz space analysis was studied by Schapira [26]. See also [4], where the case of negative multiplicities has been considered. On the other hand, to our knowledge, the $L^p$-harmonic analysis of the hypergeometric Fourier transform has not yet been developed for $p \neq 2$. The difficulty of extending the classical $L^p$-results for the spherical Fourier transform to the context of the hypergeometric Fourier transform is related to the fact that several tools coming from the geometry of the symmetric spaces are now missing. An example is Harish-Chandra’s integral formula for the spherical functions.

In this paper, our main result is Theorem 4.2, which characterizes the hypergeometric functions of Heckman and Opdam that are bounded. It is the necessary step for studying the holomorphic properties of the hypergeometric Fourier transform of $L^1$ functions. Our result is a natural extension of the celebrated theorem of Helgason and Johnson [18] characterizing the bounded spherical functions on a Riemannian symmetric space $G/K$ of the noncompact type. Suppose the triple $(a, \Sigma, m)$ arises from $G/K$. Let $\rho \in a^*$ be defined by (6), and let $C(\rho)$ denote the convex hull of the finite set $\{w\rho : w \in W\}$. The theorem of Helgason and Johnson states that the spherical function $\varphi_{\lambda}$ on $G/K$ is bounded if and only if $\lambda$ belongs to the tube domain in $a^*_C$ over $C(\rho)$. Theorem 4.2 proves that this characterization extends to Heckman-Opdam’s hypergeometric functions $\varphi_{\lambda}$ associated with any triple $(a, \Sigma, m)$. Some partial results in this direction have also been obtained by Rösler in [25], by methods that are different from ours.

Our principal tools are appropriate series expansions of Heckman and Opdam hypergeometric functions $\varphi_{\lambda}$ on a positive Weyl chamber $a^+$ of $a$.

For generic values of the spectral parameter $\lambda$, the function $\varphi_{\lambda}$ is defined on $a^+$ as linear combination of the Harish-Chandra series $\Phi_{w\lambda}$ with $w \in W$:

$$ \varphi_{\lambda} = \sum_{w \in W} c(w\lambda)\Phi_{w\lambda}, $$

where $c$ denotes Harish-Chandra’s $c$-function; see (19) and (20). By construction, the Harish-Chandra series are exponential series on $a^+$. So, for generic $\lambda$’s, one immediately obtains an exponential series expansion for $\varphi_{\lambda}$. But this expansion does not extend to all spectral parameters: its coefficients are meromorphic functions of $\lambda \in a^*_C$ and have singularities for non-generic $\lambda$. To get series expansions which hold for all values of the spectral parameter, one needs a different method.

The idea behind the method used in this paper appeared in the study of spherical functions on Riemannian symmetric spaces of the noncompact type. In this case, the spherical function $\varphi_0$ can be recovered from $\varphi_{\lambda}$ with $\lambda$ near 0. Indeed, there is a polynomial $\pi(\lambda)$ and a positive
constant $C$ so that on $a$ we have

$$\varphi_0 = C \partial(\pi)(\pi(\lambda)\varphi_\lambda)|_{\lambda=0},$$

In (3), $\partial(\pi)$ is the constant coefficient differential operator on $a$ canonically associated with $\pi$; see (22) and section 1.1 for the precise definitions of $\pi$ and $\partial(\pi)$, respectively. Formula (3) originated in the work of Harish-Chandra (see e.g. [9], p. 165 and references therein). It was applied by Anker [1] to obtain an exponential series expansion and sharp estimates for Harish-Chandra’s spherical functions $\varphi_0$. In [26], Schapira proved that the same formula holds for the Heckman-Opdam’s hypergeometric function $\varphi_0$ and extended Anker’s results to this case. More precisely, the expansion of $\varphi_0$ comes from the explicit computation of the right-hand side of (3) after substituting $\varphi_\lambda$ with its exponential series expansion (2) for generic $\lambda$’s. The point is that the polynomial $\pi(\lambda)$ cancels all singularities near $\lambda = 0$ of the meromorphic coefficients of the expansion of $\varphi_\lambda$ (and it is the minimal polynomial having this property).

Our first step is Proposition 2.5, where we prove an analog of (3) for the Heckman-Opdam’s hypergeometric functions of arbitrary spectral parameters. This requires a precise knowledge of the singularities of the coefficients of the exponential series expansion (2). For every fixed $\lambda_0 \in a^*_C$ we find a polynomial $p(\lambda)$ (depending on $\lambda_0$ and minimal in a suitable sense) so that multiplication by $p(\lambda)$ cancels all singularities of the meromorphic coefficients of the exponential series expansion (2) in a neighborhood of $\lambda_0$. Notice that the analysis needed for $\lambda_0$ arbitrary is more delicate than the one for $\lambda_0 = 0$. Indeed, near $\lambda_0 = 0$, the coefficients of the Harish-Chandra series are holomorphic. Consequently, only the singularities of Harish-Chandra’s $c$-function play a role. They are located along the root hyperplanes, hence along a Weyl-group invariant finite family of hyperplanes through 0. On the other hand, near an arbitrary point $\lambda_0$, one has to consider the singularities of the $c$-function, those of the coefficients of the Harish-Chandra series, and one also needs to understand how they transform under the action of the Weyl group. The polynomial $p(\lambda)$ is therefore a product of factors which take into account all these different contributions; see (49). The operator $\partial(\pi)$ occurring in the analog of (3) for $\lambda_0$ is the constant coefficient differential operator associated with the highest order term $\pi(\lambda)$ of $p(\lambda)$. The general version of (3),

$$\varphi_{\lambda_0} = C \partial(\pi)(p(\lambda)\varphi_\lambda)|_{\lambda=\lambda_0},$$

is thus built using two polynomials, $p(\lambda)$ and $\pi(\lambda)$, which agree in the very special case $\lambda_0 = 0$. The right-hand side of (4), together with the exponential expansion for generic $\lambda$’s, allows us to calculate the exponential series expansion of $\varphi_{\lambda_0}$; see Theorem 2.11 and Corollary 2.13. Some applications of the series expansions are then obtained in sections 3 and 4.

Besides the characterization of the set of $\lambda$’s for which the hypergeometric function $\varphi_\lambda$ is bounded, this article contains the following results in the asymptotic analysis of the Heckman-Opdam’s hypergeometric functions $\varphi_\lambda$:

1. An exponential series expansion for $\varphi_\lambda$, even when $\lambda \in a^*_C$ is not generic (Theorem 2.11).
2. Estimates for $\varphi_\lambda$ away from the walls of the Weyl chambers for all $\lambda \in a^*_C$ (Theorem 3.1).
3. Sharp estimates for $\varphi_\lambda$ when $\lambda \in a^*$ (Theorem 3.4).
The sharp estimates of Theorem 3.4 were stated without proof in [26, Remark 3.1].

In the last section of the paper we develop the $L^p$-theory for the hypergeometric Fourier transform for $0 < p < 2$. Using Theorem 4.2 we study the holomorphic properties of the hypergeometric Fourier transform on $L^p(a, du)^W$ when $1 < p < 2$. We prove the Hausdorff-Young inequalities and the Riemann-Lebesgue lemma. We also establish injectivity and an inversion formula for the hypergeometric Fourier transform. Then, by an easy generalization of Anker’s results [2] for the $L^p$-spherical Fourier transform on Riemannian symmetric spaces, we prove an $L^p$-Schwartz space isomorphism theorem for $0 < p < 2$ (see Theorem 5.6).

1. Notation and preliminaries

We shall use the standard notation $\mathbb{N}, \mathbb{N}_0, \mathbb{Z}, \mathbb{R}$ and $\mathbb{C}$ for the positive integers, the nonnegative integers, the integers, the reals and the complex numbers. The symbol $A \cup B$ denotes the union of $A$ and $B$, whereas $A \sqcup B$ indicates their disjoint union. Given two nonnegative functions $f$ and $g$ on a domain $D$, we write $f \asymp g$ if there exists positive constants $C_1$ and $C_2$ so that $C_1g(x) \leq f(x) \leq C_2g(x)$ for all $x \in D$.

Let $a$ be an $l$-dimensional real Euclidean vector space with inner product $\langle \cdot, \cdot \rangle$, and let $a^*$ be the dual space of $a$. For $\lambda \in a^*$ let $x_\lambda \in a$ be determined by $\lambda(x) = \langle x, x_\lambda \rangle$ for all $x \in a$. The assignment $\langle \lambda, \mu \rangle := \langle x_\lambda, x_\mu \rangle$ defines an inner product in $a^*$. Let $a_C$ and $a^*_C$ denote the complexifications of $a$ and $a^*$, respectively. The $\mathbb{C}$-bilinear extension to $a_C$ and $a^*_C$ of the inner products on $a^*$ and $a$ will also be denoted by $\langle \cdot, \cdot \rangle$. We shall often employ the notation

$$\lambda_\alpha := \frac{\langle \lambda, \alpha \rangle}{\langle \alpha, \alpha \rangle}.$$

We shall also set $|x| = \langle x, x \rangle^{1/2}$ for $x \in a$.

Let $\Sigma$ be a (possibly nonreduced) root system in $a^*$ with associated Weyl group $W$. For $\alpha \in \Sigma$, we denote by $r_\alpha$ the reflection $\lambda \mapsto \lambda - 2\lambda_\alpha \alpha$ in $a^*$. For a set $\Sigma^+ \subseteq \Sigma$ of positive roots in $\Sigma$, let $\Pi = \{\alpha_1, \ldots, \alpha_l\} \subset \Sigma^+$ denote the corresponding set of simple roots. We denote by $\Sigma_0$ the indivisible roots in $\Sigma$: if $\alpha \in \Sigma_0$, then $\alpha/2 \notin \Sigma$. We set $\Sigma_0^+ = \Sigma^+ \cap \Sigma_0$.

A positive multiplicity function on $\Sigma$ is a $W$-invariant function $m : \Sigma \to ]0, +\infty[$. Setting $m_\alpha := m(\alpha)$ for $\alpha \in \Sigma$, we therefore have $m_{w\alpha} = m_\alpha$ for all $w \in W$. We extend $m$ to $a^*$ by putting $m_\alpha = 0$ for $\alpha \notin \Sigma$. We say that a multiplicity function $m$ is geometric if there is a Riemannian symmetric space of noncompact type $G/K$ with restricted root system $\Sigma$ such that $m_\alpha$ is the multiplicity of the root $\alpha$ for all $\alpha \in \Sigma$. Otherwise, $m$ is said to be non-geometric.

The dimension $l$ of $a$ will also be called the (real) rank of the triple $(a, \Sigma, m)$.

In this paper we adopt the notation commonly used in the theory of symmetric spaces. It differs from the notation in the work of Heckman and Opdam in the following ways. The root system $R$ and the multiplicity function $k$ used by Heckman and Opdam are related to our $\Sigma$ and $m$ by the relations $R = \{2\alpha : \alpha \in \Sigma\}$ and $k_\alpha = m_\alpha/2$ for $\alpha \in \Sigma$.

We view $a_C$ of $a$ as the Lie algebra of the complex torus $A_C := a_C/\mathbb{Z}\{2\pi i \alpha_\alpha/\langle \alpha, \alpha \rangle : \alpha \in \Sigma\}$. We write $\exp : a_C \to A_C$ for the exponential map, with multi-valued inverse log. The split real form $A := \exp a$ of $A_C$ is an abelian subgroup with Lie algebra $a$ such that $\exp : a \to A$ is a diffeomorphism. In the following, to simplify the notation, we shall identify $A$ with $a$ by means of this diffeomorphism.
The action of $W$ extends to $\mathfrak{a}$ by duality, to $\mathfrak{a}_C$ and $\mathfrak{a}_C$ by $\mathbb{C}$-linearity. Moreover, $W$ acts on functions $f$ on any of these spaces by $(wf)(x) := f(w^{-1}x)$, $w \in W$.

The positive Weyl chamber $\mathfrak{a}^+$ consists of the elements $x \in \mathfrak{a}$ for which $\alpha(x) > 0$ for all $\alpha \in \Sigma^+$; its closure is $\overline{\mathfrak{a}^+} = \{ x \in \mathfrak{a} : \alpha(x) \geq 0 \text{ for all } \alpha \in \Sigma^+ \}$. Dually, the positive Weyl chamber $(\mathfrak{a}^*)^+$ consists of the elements $\lambda \in \mathfrak{a}^*$ for which $\langle \lambda, \alpha \rangle > 0$ for all $\alpha \in \Sigma^+$. Its closure is denoted $(\mathfrak{a}^*)^+$. We write $\lambda \leq \mu$ if $\lambda, \mu \in \mathfrak{a}^*$ and $\mu - \lambda \in (\mathfrak{a}^*)^+$. The sets $\overline{\mathfrak{a}^+}$ and $(\mathfrak{a}^*)^+$ are fundamental domains for the action of $W$ on $\mathfrak{a}$ and $\mathfrak{a}^*$, respectively.

The restricted weight lattice of $\Sigma$ is $P = \{ \lambda \in \mathfrak{a}^* : \lambda_\alpha \in \mathbb{Z} \text{ for all } \alpha \in \Sigma \}$. Observe that $\{ 2\alpha : \alpha \in \Sigma \} \subseteq P$. If $\lambda \in P$, then the exponential $e^\lambda : A_C \to \mathbb{C}$ given by $e^\lambda(h) := e^{\lambda(\log h)}$ is single valued. The $e^A$ are the algebraic characters of $A_C$. Their $\mathbb{C}$-linear span coincides with the ring of regular functions $\mathbb{C}[A_C]$ on the affine algebraic variety $A_C$. The lattice $P$ is $W$-invariant, and the Weyl group acts on $\mathbb{C}[A_C]$ according to $w(e^\lambda) := e^{w\lambda}$. The set $A_C^{reg} := \{ h \in A_C : e^{2\alpha(\log h)} \neq 1 \text{ for all } \alpha \in \Sigma \}$ consists of the regular points of $A_C$ for the action of $W$. Notice that $A^+ \equiv \mathfrak{a}^+$ is a subset of $A_C^{reg}$. The algebra $\mathbb{C}[A_C^{reg}]$ of regular functions on $A_C^{reg}$ is the subalgebra of the quotient field of $\mathbb{C}[A_C]$ generated by $\mathbb{C}[A_C]$ and by $1/(1 - e^{-2\alpha})$ for $\alpha \in \Sigma^+$. Its $W$-invariant elements form the subalgebra $\mathbb{C}[A_C^{reg}]^W$.

1.1. Cherednik operators and the hypergeometric system. In this subsection we outline the theory of hypergeometric differential equations associated with root systems. This theory has been developed by Heckman, Opdam and Cherednik. We refer the reader to [3], [13], [15], [22], [23] for more details and further references.

Let $S(\mathfrak{a}_C)$ denote the symmetric algebra over $\mathfrak{a}_C$ considered as the space of polynomial functions on $\mathfrak{a}_C^*$, and let $S(\mathfrak{a}_C)^W$ be the subalgebra of $W$-invariant elements. Every $p \in S(\mathfrak{a}_C)$ defines a constant-coefficient differential operators $\partial(p)$ on $A_C$ and on $\mathfrak{a}_C$ such that $\partial(x) = \partial_x$ is the directional derivative in the direction of $x$ for all $x \in \mathfrak{a}$. The algebra of the differential operators $\partial(p)$ with $p \in S(\mathfrak{a}_C)$ will also be indicated by $S(\mathfrak{a}_C)$. Let $\mathbb{D}(A_C^{reg}) := \mathbb{C}[A_C^{reg}] \otimes S(\mathfrak{a}_C)$ denote the algebra of differential operators on $A_C$ with coefficients in $\mathbb{C}[A_C^{reg}]$. The Weyl group $W$ acts on $\mathbb{D}(A_C^{reg})$ according to

$$w(\varphi \otimes \partial(p)) := w\varphi \otimes \partial(wp).$$

We write $\mathbb{D}(A_C^{reg})^W$ for the subspace of $W$-invariant elements. The space $\mathbb{D}(A_C^{reg}) \otimes \mathbb{C}[W]$ can be endowed with the structure of an associative algebra with respect to the product

$$(D_1 \otimes w_1) \cdot (D_2 \otimes w_2) = D_1 w_1(D_2) \otimes w_1 w_2,$$

where the action of $W$ on differential operators is defined by $(wD)(wf) := w(Df)$ for every sufficiently differentiable function $f$. It is also a left $\mathbb{C}[A_C^{reg}]$-module. Considering $D \in \mathbb{D}(A_C^{reg})$ as an element of $\mathbb{D}(A_C^{reg}) \otimes \mathbb{C}[W]$, we shall usually write $D$ instead of $D \otimes 1$. The elements of the algebra $\mathbb{D}(A_C^{reg}) \otimes \mathbb{C}[W]$ are called the differential-reflection operators on $A_C^{reg}$. The differential-reflection operators act on functions $f$ on $A_C^{reg}$ according to $(D \otimes w)f := D(wf)$.

For $x \in \mathfrak{a}$ the Cherednik operator (or Dunkl-Cherednik operator) $T_x \in \mathbb{D}(A_C^{reg}) \otimes \mathbb{C}[W]$ is defined by

$$T_x := \partial_x - \rho(x) + \sum_{\alpha \in \Sigma^+} m_\alpha \alpha(x)(1 - e^{-2\alpha})^{-1} \otimes (1 - r_\alpha).$$
\[ \rho := \frac{1}{2} \sum_{\alpha \in \Sigma^+} m_\alpha \alpha \in \mathfrak{a}^*. \]

The Cherednik operators can also be considered as operators acting on smooth functions on \( \mathfrak{a} \). This is possible because, as can be seen from the Taylor formula, the term \( 1 - r_\alpha \) cancels the apparent singularity arising from the denominator \( 1 - e^{-2\alpha} \). The Cherednik operators commute with each other; cf. [22, Section 2]. Therefore the map \( x \mapsto T_x \) on \( \mathfrak{a} \) extends uniquely to an algebra homomorphism \( p \mapsto T_p \) of \( S(\mathfrak{a}_C) \) into \( \mathcal{D}(A_{C}^{\text{reg}}) \otimes \mathbb{C}[W] \).

Define a linear map \( \Upsilon : \mathcal{D}(A_{C}^{\text{reg}}) \otimes \mathbb{C}[W] \to \mathcal{D}(A_{C}^{\text{reg}}) \) by
\[
\Upsilon(\sum_j D_j \otimes w_j) := \sum_j D_j.
\]

Then \( \Upsilon(Q)f = Qf \) for all \( Q \in \mathcal{D}(A_{C}^{\text{reg}}) \otimes \mathbb{C}[W] \) and all \( W \)-invariant \( f \) on \( A_{C}^{\text{reg}} \).

For \( p \in S(\mathfrak{a}_C) \) we set \( D_p := \Upsilon(T_p) \). If \( p \in S(\mathfrak{a}_C)^W \), then \( D_p \in \mathcal{D}(A_{C}^{\text{reg}})^W \); see [22, Theorem 2.12(2)]. The algebra
\[ \mathbb{D} = \mathbb{D}(\mathfrak{a}, \Sigma, m) := \{ D_p : p \in S(\mathfrak{a}_C)^W \} \]
is a commutative subalgebra of \( \mathcal{D}(A_{C}^{\text{reg}})^W \). It is called the algebra of hypergeometric differential operators associated with \( (\mathfrak{a}, \Sigma, m) \). It is the analogue, for arbitrary multiplicity functions, of the commutative algebra of the radial components on \( A = \exp \mathfrak{a} \) of the invariant differential operators on a Riemannian symmetric space of noncompact type.

A remarkable element of \( \mathbb{D} \) corresponds to the polynomial \( p_L \in S(\mathfrak{a}_C)^W \) defined by \( p_L(\lambda) := \langle \lambda, \lambda \rangle \) for \( \lambda \in \mathfrak{a}_C^* \). Then
\[ D_{p_L} = L + \langle \rho, \rho \rangle, \]
where
\[ L := L_\alpha + \sum_{\alpha \in \Sigma^+} m_\alpha \coth \alpha \partial_\alpha \]
and \( L_\alpha \) is the Laplace operator on \( \mathfrak{a} \); see [13, Theorem 2.2]. In (7) we have set \( \partial_\alpha := \partial(x_\alpha) \) and
\[ \coth \alpha := \frac{1 + e^{-2\alpha}}{1 - e^{-2\alpha}}. \]

If \( (\mathfrak{a}, \Sigma, m) \) is geometric, then \( L \) coincides with the radial component on \( \mathfrak{a} \equiv A \) with respect to the left action of \( K \) of the Laplace operator on a Riemannian symmetric space \( G/K \) of noncompact type.

The map \( \gamma : \mathbb{D} \to S(\mathfrak{a}_C)^W \) defined by
\[ \gamma(D_p)(\lambda) := p(\lambda) \]
is called the Harish-Chandra homomorphism. It defines an algebra isomorphism of \( \mathbb{D} \) onto \( S(\mathfrak{a}_C)^W \) (see [15, Theorem 1.3.12 and Remark 1.3.14]). From Chevalley’s theorem it therefore follows that \( \mathbb{D} \) is generated by \( l (= \dim \mathfrak{a}) \) elements.

Let \( \lambda \in \mathfrak{a}_C^* \) be fixed. The system of differential equations
\[ D_p \varphi = p(\lambda) \varphi, \quad p \in S(\mathfrak{a}_C)^W, \]
is called the hypergeometric differential equations.
is called the hypergeometric system of differential equations with spectral parameter $\lambda$ associated with the data $(a, \Sigma, m)$. The differential equation corresponding to the polynomial $p_L$ is

$$L\varphi = (\langle \lambda, \lambda \rangle - \langle \rho, \rho \rangle)\varphi.$$  

For geometric multiplicities, the hypergeometric system (9) agrees with the system of differential equations defining Harish-Chandra’s spherical function of spectral parameter $\lambda$.

**Example 1.1** (The rank-one case). The rank-one case corresponds to triples $(a, \Sigma, m)$ in which $a$ is one dimensional. Then the set $\Sigma^+$ consists at most of two elements: $a$ and, possibly, $2a$. By setting $x_\alpha/2 \equiv 1$ and $\alpha \equiv 1$, we identify $a$ and $a^\circ$ with $\mathbb{R}$, and their complexifications $a_\mathbb{C}$ and $a_\mathbb{C}^\circ$ with $\mathbb{C}$. The Weyl chamber $a^+$ coincides with the half-line $]0, +\infty[. The Weyl group $W$ reduces to $\{-1, 1\}$ acting on $\mathbb{R}$ and $\mathbb{C}$ by multiplication. The algebra $\mathbb{D}$ is generated by $D_{p_L} = L + \rho^2$. The hypergeometric differential system with spectral parameter $\lambda \in \mathbb{C}$ is equivalent to the single Jacobi differential equation

$$d^2\varphi \over dz^2 + (m_\alpha \coth z + m_{2\alpha} \coth(2z)) \, d\varphi \over dz = (\lambda^2 - \rho^2)\varphi.$$  

The function $z \mapsto e^z$ maps $a_\mathbb{C} \equiv \mathbb{C}$ onto $A_\mathbb{C} \equiv \mathbb{C}^\times$. Hence $A_\mathbb{C}^{\text{reg}} \equiv \mathbb{C} \setminus \{0, \pm 1\}$. The change of variable $\zeta := (1 - \cosh z)/2$ transforms (11) into the hypergeometric differential equation

$$\zeta(1 - \zeta) {d^2\psi \over d\zeta^2} + [c - (1 + a + b)\zeta] {d\psi \over d\zeta} - ab\zeta = 0$$  

with parameters

$$a = \frac{\lambda + \rho}{2}, \quad b = \frac{-\lambda + \rho}{2}, \quad c = \frac{m_\alpha + m_{2\alpha} + 1}{2}.$$  

1.2. The Harish-Chandra series. As in the classical theory of spherical functions on Riemannian symmetric spaces, the explicit expression of the differential equation (10) suggested to Heckman and Opdam [14] to look for solutions on $a^+$ of the hypergeometric system (9) with spectral parameter $\lambda$ which are of the form

$$\Phi_\lambda(x) = e^{(\lambda - \rho)(x)} \sum_{\mu \in \Lambda} \Gamma_\mu(\lambda) e^{-\mu(x)}, \quad x \in a^+.$$  

Here $\Lambda := \left\{ \sum_{j=1}^l n_j \alpha_j : n_j \in \mathbb{N}_0 \right\}$ is the positive semigroup generated by the fundamental system of simple roots $\Pi := \{\alpha_1, \ldots, \alpha_l\}$ in $\Sigma^+$. For $\mu \in \Lambda \setminus \{0\}$, the coefficients $\Gamma_\mu(\lambda)$ are rational functions of $\lambda \in a_\mathbb{C}^*$ determined from the recursion relations

$$\langle \mu, \mu - 2\lambda \rangle \Gamma_\mu(\lambda) = 2 \sum_{\alpha \in \Sigma^+} m_\alpha \sum_{\kappa \in \mathbb{N}} \Gamma_{\mu - 2k\alpha}(\lambda) \langle \mu + \rho - 2k\alpha - \lambda, \alpha \rangle,$$  

with initial condition $\Gamma_0(\lambda) = 1$. They are derived by formally inserting the series for $\Phi$ into the differential equation (10). Let $\ell(\mu) := \sum_{j=1}^l n_j$ denote the level of $\mu = \sum_{j=1}^l n_j \alpha_j \in \Lambda$. It is easy to check by induction on $\ell(\mu)$ that the recursion relations imply $\Gamma_\mu(\lambda) = 0$ unless $\mu = \sum_{j=1}^l n_j \alpha_j$ with $n_j \geq 0$ and $n_j$ even for all $j = 1, \ldots, l$. Hence the function $\Phi_\lambda(x)$ is in
fact a sum over \(2\Lambda\), that is
\[
\Phi_\lambda(x) = e^{(\lambda - \rho)(x)} \sum_{\mu \in 2\Lambda} \Gamma_\mu(x) e^{-\mu(x)}, \quad x \in a^+.
\]

The function \(\Phi_\lambda(x)\) is called the \textit{Harish-Chandra series}.

Let \(\mu \in 2\Lambda \setminus \{0\}\) be fixed. A priori, the relation \([12]\) uniquely define the rational function \(\Gamma_\mu(\lambda)\) on \(a^*_C\) provided \(\langle \tau, \tau - 2\lambda \rangle \neq 0\) for all \(\tau \in 2\Lambda \setminus \{0\}\) with \(\tau \leq \mu\). Opdam proved that, in fact, many of these singularities are removable. Correspondingly, many of the apparent singularities of the Harish-Chandra series are removable as well.

In the following we adopt the notation
\[
H_{\alpha,r} = \{ \lambda \in a^*_C : \lambda_\alpha = r \}.
\]

We shall consider meromorphic functions \(f\) on \(a^*_C\) with singularities on a locally finite (generically infinite) family \(H\) of affine complex hyperplanes \(H_{\alpha,r}\). We say that \(f\) has at most a \textit{simple pole along} \(H_{\alpha,r}\) if the function \(\lambda \mapsto (\lambda_\alpha - r) f(\lambda)\) extends holomorphically to a neighborhood of \(H_{\alpha,r} \setminus \bigcup_{H \in H, H \neq H_{\alpha,r}} H\).

**Theorem 1.2.** (a) Let \(\mu \in 2\Lambda \setminus \{0\}\). Then the rational function \(\Gamma_\mu(\lambda)\) has at most simple poles located along the hyperplanes \(H_{\alpha,n}\) with \(\alpha \in \Sigma_0^+, n \in \mathbb{N}\) and \(2\alpha \leq \mu\).

(b) There is a tubular neighborhood \(U^+\) of \(A^+ = \exp a^+\) in \(A_C\) so that the Harish-Chandra series \(\Phi_\lambda(x)\) is a meromorphic function of \((\lambda, x) \in a^*_C \times U^+\) with at most simple poles along hyperplanes of the form \(H_{\alpha,n}\) with \(\alpha \in \Sigma_0^+\) and \(n \in \mathbb{N}\).

**Proof.** This is Corollary 2.10 in \([21]\). See also \([15\), Proposition 4.2.5\] and \([23\), Lemma 6.5\]. \(\square\)

**Remark 1.3.** The neighborhood \(U^+\) in Theorem \([12]\) can be chosen of the form \(A^+ U_0\) where \(U_0\) is a connected and simply connected neighborhood of \(e = \exp 0\) in \(T = \exp(i a)\) so that the function \(\log\) is single valued on it. Then all functions \(e^{(\lambda - \rho)(\log h)}\) (with \(\lambda \in a^*_C\) and \(h \in A_C\)) are single valued and holomorphic on \(A^+ U_0\).

The convergence of the Harish-Chandra series can be studied by estimating its coefficients. We record the following result, which is due to Opdam (\([21\), Lemma 2.1]; see also \([15\), Lemma 4.4.2\]). It is an extension of the classical argument by Helgason in \([16\), Lemma 4.1]; see also \([17\), Ch. IV, Lemma 5.3\). We state it in a slightly modified form (fixed multiplicity function and variable Weyl group element), which is more suitable to our purposes. The last part of the lemma is a consequence of the first part and Cauchy’s integral formula.

**Lemma 1.4.** Let \(U \subset a^*_C\) be an open set with compact closure \(\overline{U}\), and let \(w \in W\). Let \(d(\lambda)\) be a holomorphic function such that \(d(\lambda) \Gamma_\mu(w \lambda)\) is holomorphic on \(\overline{U}\) and all \(\mu \in 2\Lambda \setminus \{0\}\).

Let \(x_0 \in a^+\) be fixed. Then there is a constant \(M_{U,x_0}\) such that
\[
|d(\lambda) \Gamma_\mu(w \lambda)| \leq M_{U,x_0} e^{\mu(x_0)}
\]
for all \(\mu \in 2\Lambda\) and \(\lambda \in U\). Hence the series
\[
e^{(w - \rho)(x)} \sum_{\mu \in 2\Lambda} d(\lambda) \Gamma_\mu(w \lambda) e^{-\mu(x)}
\]
converges absolutely and uniformly in \((\lambda, x) \in U \times (x_0 + \overline{a^+})\) to \(d(\lambda) \Phi_{w \lambda}(x)\). Furthermore, for every \(p \in S(a_C)\) there is a constant \(M_{p,U,x_0}\) such that
\[
\left| \partial(p) \left( d(\lambda) \Gamma_\mu(w \lambda) \right) \right| \leq M_{p,U,x_0} e^{\mu(x_0)}
\]
for all \( \mu \in 2\Lambda \) and \( \lambda \in U \). The series (15) can therefore be differentiated term-by-term and the differentiated series converges absolutely and uniformly in \((\lambda, x) \in U \times (x_0 + a^+)\) to \(\partial(p)(d(\lambda)\Phi_{w\lambda}(x))\).

Notice that the explicit expression of the holomorphic function \(d\) is not relevant. We choose a specific function in Proposition 2.5.

As in the Riemannian case, the Harish-Chandra series can be used to build a basis for the smooth solutions on \(a^+\) of the entire hypergeometric system with spectral parameter \(\lambda\). This is possible when \(\lambda \in a^*_C\) is generic.

**Definition 1.5.** We say that \(\lambda \in a^*_C\) is generic if \(\lambda_\alpha \notin \mathbb{Z}\) for all \(\alpha \in \Sigma\).

Notice that, since \(\lambda_\alpha = 2\lambda_\alpha\), the element \(\lambda \in a^*_C\) is generic if and only if \(\lambda_\alpha \notin \mathbb{Z}\) for all \(\alpha \in \Sigma\). Moreover, since \((r_\beta \lambda)_\alpha = \lambda_{r_\beta \alpha}\) for all \(\alpha, \beta \in \Sigma\), the set of generic elements in \(a^*_C\) is \(W\)-invariant.

**Theorem 1.6.** Let \(U^+\) be the tubular neighborhood of \(A^+\) from Theorem 1.2. If \(\lambda \in a^*_C\) is generic, then the set \(\{\Phi_{w\lambda}(x) : w \in W\}\) is a basis of the solution space on \(U^+\) of the hypergeometric system (9) with spectral parameter \(\lambda\).

**Proof.** See [13, Corollary 4.2.6] or [23, Theorem 6.7].

**Example 1.7 (The rank-one case).** The solution of the Jacobi differential equation (11) on \((0, +\infty)\) that behaves asymptotically as \(e^{(\lambda-\rho)t}\) for \(t \to +\infty\) is

\[
\Phi_\lambda(t) = (2 \sinh t)^{\lambda-\rho} \mathcal{F}_1 \left( \frac{\rho-\lambda}{2}, \frac{-m_\alpha/2 + 1 - \lambda}{2}; 1 - \lambda; -\sinh^2 t \right),
\]

where \(\mathcal{F}_1\) denotes the Gaussian hypergeometric function; see e.g. [6, Chapter 2]. The function \(\Phi_\lambda(t)\) coincides with the Jacobi function of second kind \(\Phi^{(a,b)}_\nu(t)\) with parameters \(a = (m_\alpha + m_2\alpha - 1)/2, b = (m_2\alpha - 1)/2\) and \(\nu = -i\lambda\) (see [20, Section 2]).

**Example 1.8 (The complex case).** Let \(m\) be geometric multiplicity of a root system \(\Sigma\) which is reduced (i.e. \(2\alpha \notin \Sigma\) for every \(\alpha \in \Sigma\)). If \(m_\alpha = 2\) for all \(\alpha \in \Sigma\), then \(m\) corresponds to a Riemannian symmetric space of the noncompact type \(G/K\) with \(G\) complex. The triple \((a, \Sigma, m)\) will be said to correspond to a complex case. In the complex case, we have

\[
\Phi_\lambda(x) = \Delta(x)^{-1} e^{\lambda(x)}.
\]

where

\[
\Delta := \prod_{\alpha \in \Sigma^+} (e^\alpha - e^{-\alpha})
\]

is the Weyl denominator.

1.3. **The c-function.** For \(\alpha \in \Sigma^+_0\) and \(\lambda \in a^*_C\) we set

\[
c_\alpha(\lambda) = \frac{2^{-\lambda_\alpha}}{\Gamma\left(\frac{\lambda_\alpha}{2} + \frac{m_\alpha}{4} + \frac{1}{2}\right) \Gamma\left(\frac{\lambda_\alpha}{2} + \frac{m_\alpha}{4} + \frac{m_2\alpha}{4}\right)}.
\]
where $\Gamma$ is the Euler gamma function. Harish-Chandra’s $c$-function is the meromorphic function on $a_C^*$ defined by
\[
c(\lambda) = c_{HC} \prod_{\alpha \in \Sigma_0^+} c_{\alpha}(\lambda)
\]
where $c_{HC}$ is a normalizing constant chosen so that $c(\rho) = 1$.

Recall the notation $H_{\alpha,n} := \{ \lambda \in a_C^* : \lambda_\alpha = r \}$ from (14). Observe that the equality $H_{\alpha,n} = H_{\beta,m}$ with $\alpha, \beta \in \Sigma_0^+$ and $n, m \in \mathbb{Z}$ implies $\alpha = \beta$ and $n = m$. From the singularities of the gamma function we therefore obtain the following lemma.

**Lemma 1.9.** The meromorphic function $c(\lambda)$ admits at most simple poles located along the hyperplanes $H_{\alpha,-n}$ with $\alpha \in \Sigma_0^+$ and $n \in \mathbb{N}_0$.

The possible zeros of $c(\lambda)$ are located along the hyperplanes
\[
H_{\alpha,-(m_\alpha/2+m_2)\alpha^{-2n}} \quad \text{with} \quad \alpha \in \Sigma_0^+ \quad \text{and} \quad n \in \mathbb{N}_0,
\]
\[
H_{\alpha,-m_\alpha/2-1-2n} \quad \text{with} \quad \alpha \in \Sigma_0^+ \quad \text{and} \quad n \in \mathbb{N}_0.
\]

1.4. The hypergeometric functions of Heckman and Opdam. Let $\lambda \in a_C^*$. The hypergeometric function of spectral parameter $\lambda$ is the unique analytic $W$-invariant function $\varphi_\lambda(x)$ on $a$ which satisfies the system of differential equations (9) and which is normalized by $\varphi_\lambda(0) = 1$. In the geometric case, with the identification of $a$ with $A = \exp a$, the function $\varphi_\lambda$ agrees with the (elementary) spherical function of spectral parameter $\lambda$. For $x \in a^+$ and generic $\lambda \in a_C^*$, the hypergeometric function admits the representation
\[
\varphi_\lambda(x) = \sum_{w \in W} c(w\lambda)\Phi_w^\lambda(x).
\]

**Example 1.10 (The rank-one case).** In the rank-one case, with the identifications introduced in Example 1.1, Heckman-Opdam’s hypergeometric function coincides with the Jacobi function of the first kind
\[
\varphi_\lambda(t) = \binom{m_\alpha/2 + m_2 + \lambda}{2} \binom{m_\alpha/2 + m_2 - \lambda}{2} \binom{m_\alpha + m_2 + 1}{2} \frac{-\sinh^2 t}{\Delta(x)}
\]

**Example 1.11 (The complex case).** In the complex case the multiplicity is geometric. The hypergeometric functions of Heckman and Opdam agree with Harish-Chandra’s spherical functions. In this very special case, they are given by the explicit formula
\[
\varphi_\lambda(x) = \frac{\pi(\rho)}{\pi(\lambda)} \frac{\sum_{w \in W} (\det w) e^{w\lambda(x)}}{\Delta(x)}
\]

where
\[
\pi(\lambda) = \prod_{\alpha \in \Sigma_0^+} \langle \lambda, \alpha \rangle
\]
and $\Delta$ is as in (17). See e.g. [9, p. 251].

The nonsymmetric hypergeometric function of spectral parameter $\lambda$ is the unique analytic function $G_\lambda(x)$ on $a$ which satisfies the system of differential-reflection equations
\[
T_x G_\lambda = \lambda(x) G_\lambda, \quad x \in a,
\]
and which is normalized by $G_\lambda(0) = 1$. We have the relation

$$\varphi_\lambda(x) = \frac{1}{|W|} \sum_{w \in W} G_\lambda(wx), \quad x \in \mathfrak{a}.$$  \hspace{1cm} (24)

Schapira proved in [26] that $\varphi_\lambda$ is real and strictly positive for $\lambda \in \mathfrak{a}^*$. He also proved the fundamental estimate:

$$|\varphi_\lambda| \leq \varphi_{\Re \lambda}, \quad \lambda \in \mathfrak{a}^*_+.$$  \hspace{1cm} (25)

We refer to [22], [15] and [26] for the proof of these statements and for further information.

1.5. The hypergeometric Fourier transform. Let $dx$ denote a fixed normalization of the Haar measure on $\mathfrak{a}$. We associate with the triple $(\mathfrak{a}, \Sigma, m)$ the measure $d\mu(x) = \mu(x) \, dx$ on $\mathfrak{a}$, where

$$\mu(x) = \prod_{\alpha \in \Sigma^+} |e^{\alpha(x)} - e^{-\alpha(x)}|^{m_\alpha}.$$  \hspace{1cm} (26)

Notice that when $(\mathfrak{a}, \Sigma, m)$ comes from a Riemannian symmetric space $G/K$, then $d\mu$ is the component along $A \equiv \mathfrak{a}$ of the Haar measure on $G$ with respect to the Cartan decomposition $G = KAK$.

Recall from (1) the definition of the hypergeometric Fourier transform $\mathcal{F}f = \hat{f}$ of a sufficiently regular $W$-invariant functions on $\mathfrak{a}$.

Let $\Gamma$ be a $W$-invariant compact convex subset of $\mathfrak{a}$ and let $q_\Gamma(\lambda) = \sup_{x \in \Gamma} \lambda(x)$, with $\lambda \in \mathfrak{a}$, be the supporting function of $\Gamma$. The Paley-Wiener space $PW_\Gamma(\mathfrak{a}_c^*)^W$ consists of all $W$-invariant entire functions $F$ on $\mathfrak{a}_c^*$ satisfying

$$\sup_{\lambda \in \mathfrak{a}_c^*} (1 + |\lambda|)^N e^{-q_\Gamma(\Re \lambda)} |F(\lambda)| < \infty$$  \hspace{1cm} (27)

for all $N \in \mathbb{N}_0$. We topologize $PW_\Gamma(\mathfrak{a}_c^*)^W$ by the seminorms defined by the left-hand side of (27). Let $C^\infty_\Gamma(\mathfrak{a})^W$ denote the space of $W$-invariant smooth functions on $\mathfrak{a}$ with support inside $\Gamma$. The space $C^\infty_\Gamma(\mathfrak{a})^W$ is considered with the topology induced by $C^\infty_c(\mathfrak{a})$. We will only be interested in two situations: when $\Gamma$ is equal to $B_R := \{x \in \mathfrak{a} : |x| \leq R\}$ for some $R > 0$, and when $\Gamma$ is equal to the polar set $C_\Lambda$ of the convex hull of the set $\{w(\Lambda) : w \in W\}$ (with $\Lambda \in \mathfrak{a}^*$). Recall that $C_\Lambda = \{x \in \mathfrak{a} : \Lambda(x^+) \leq 1\}$, where $x^+$ is the unique element in $\mathfrak{a}^+$ of the Weyl group orbit of $x$. In the first case, $q_{B_R}(\lambda) = R|\lambda|$. Hence, the usual Paley-Wiener space on $\mathfrak{a}_c^*$ denoted by $PW(\mathfrak{a}_c^*)^W$, is the union of the spaces $PW_{B_R}(\mathfrak{a}_c^*)^W$, with the inductive limit topology.

The following theorem gives the basic results in the $L^2$-harmonic analysis of the hypergeometric Fourier transform. It is due to Opdam [22], who proved the Paley-Wiener theorem in a slightly less general form. In the classical geometric case of Riemannian symmetric spaces of the noncompact type and with $\Gamma = B_R$, the theorem is due to Helgason [16] and Gangolli [8]; with $\Gamma = C_\Lambda$ it was proven by Anker [2]. The fact that Anker’s results also extend to the non-geometric case was observed by Schapira in [26, p. 240].

**Theorem 1.12.** We keep the above notation and assumptions.

(a) (Paley-Wiener theorem) The hypergeometric Fourier transform $\mathcal{F}$ is a topological isomorphism between $C^\infty_c(\mathfrak{a})^W$ and $PW(\mathfrak{a}_c^*)^W$. It restricts to a topological isomorphism between $C^\infty_\Gamma(\mathfrak{a})^W$ and $PW_\Gamma(\mathfrak{a}_c^*)^W$. 


Lemma 2.1. Let \( \lambda \in a_+^* \) under the assumption that \( x \in a^+ \) is sufficiently far from the walls of \( a^+ \). We begin with some properties of the centralizer of \( \lambda \) that will be needed in the sequel.

We shall employ the notation \( W_\lambda = \{ w \in W : w\lambda = \lambda \} \) for the centralizer of \( \lambda \in a_+^* \) in \( W \). For \( \Theta \subset \Pi \) we denote by \( W_\Theta \) the (standard parabolic) subgroup of \( W \) generated by the reflections \( r_\alpha \) with \( \alpha \in \Theta \). Moreover, we write \( \Sigma_\Theta \) for the subsystem of \( \Sigma \) consisting of the roots which can be written as linear combinations of elements from \( \Theta \). Furthermore, we set \( \Sigma_\Theta^+ = \Sigma_\Theta \cap \Sigma^+ \) and \( \Sigma_{\Theta,0}^+ = \Sigma_\Theta \cap \Sigma_0^+ \).

Recall the notation \( \lambda_\alpha \) from [3]. For \( \lambda \in a_+^* \) we define

\[
\lambda_0 = \{ \alpha \in \Sigma_0^+ : \lambda_\alpha \in \mathbb{Z} \} ,
\]

\[
\lambda_\alpha = \{ \alpha \in \Sigma_0^+ : \lambda_\alpha \in \mathbb{N} \} ,
\]

\[
\lambda_\alpha = \{ \alpha \in \Sigma_0^+ : \lambda_\alpha = 0 \} .
\]

Suppose that \( \text{Re} \lambda \in (a^*)^+ \). Then \( W_{\text{Re} \lambda} = W_{\Theta(\lambda)} \) where \( \Theta(\lambda) = \Sigma_0^+ \cap \Pi \). Observe that \( W_{\Theta(\lambda)} \) is the Weyl group of \( \Sigma_{\Theta(\lambda)} \) and \( \Sigma_{\Theta(\lambda),0}^+ = \Sigma_0^+ \).

Let \( w_\lambda \in W \) be chosen so that \( w_\lambda \text{Im} \lambda \in (a^*)^+ \). Then there is \( \Xi(\lambda) \subset \Pi \) so that \( W_{\text{Im} \lambda} = w_\lambda^{-1}W_{\Xi(\lambda)}w_\lambda = W_{w_\lambda^{-1}\Xi(\lambda)} \). Moreover, \( \Sigma_{w_\lambda^{-1}\Xi(\lambda),0}^+ = \Sigma_0 \).

**Lemma 2.1.** Let \( \lambda \in a_+^* \) with \( \text{Re} \lambda \in (a^*)^+ \), and keep the above notation. Then

\[
W_\lambda = W_{\text{Re} \lambda} \cap W_{\text{Im} \lambda} = W_{\Theta(\lambda)} \cap W_{w_\lambda^{-1}\Xi(\lambda)} ,
\]

\[
\Sigma_0 = \Sigma_{\Theta(\lambda),0}^+ \cap \Sigma_{w_\lambda^{-1}\Xi(\lambda),0}^+ \subset \Sigma_{\Theta(\lambda),0}^+ ,
\]

\[
\Sigma_\lambda = \Sigma_0^+ \cup \Sigma_\lambda^+ ,
\]

\[
\Sigma_\lambda^+ \cap \Sigma_{\Theta(\lambda),0}^+ = \emptyset .
\]

Moreover, let \( w \in W_\lambda \). Then

\[
w(\Sigma_0^+) \subset \Sigma_0^+ \backslash \Sigma_\lambda^+ ,
\]

\[
w(\Sigma_0^+ \backslash \Sigma_0^+) = \Sigma_0^+ \backslash \Sigma_\lambda^+ .
\]

Furthermore, if \( w \in W_{\text{Re} \lambda} = W_{\Theta(\lambda)} \), then

\[
w(-\Sigma_0^+) \cap \Sigma_\lambda^+ = \emptyset ,
\]

\[
w(\Sigma_0^+) \cap \Sigma_\lambda^+ = \Sigma_\lambda^+ .
\]
Proof. The first part of the Lemma is an immediate consequence of the discussion above and the fact that $\text{Re } \lambda \in (\mathfrak{a}^*)^\perp$. For (34), observe that if $\alpha \in \Sigma_\lambda^+$ then $\langle \text{Re } \lambda, \alpha \rangle = \langle \lambda, \alpha \rangle > 0$. To prove (35), notice that $w^{-1} \in W_\lambda$. Hence for $\alpha \in \Sigma_\lambda^0$ we have $\lambda_{\alpha\alpha} = (w^{-1})\alpha = \lambda_\alpha = 0$. Formula (36) is a consequence of (35) applied to $w$ and $w^{-1}$ and the fact that $W_\lambda$ is the Weyl group of the closed subsystem $\Sigma_\lambda^0 \sqcup (-\Sigma_\lambda^0)$ of $\Sigma$.

Suppose now that $w \in W_{\Theta(\lambda)}$, so $w(\text{Re } \lambda) = \text{Re } \lambda$. Let $\alpha \in w(-\Sigma_\lambda^+)$. Then

$$\text{Re}(\lambda_\alpha) = (\text{Re } \lambda)_\alpha = (w(\text{Re } \lambda))_\alpha = (\text{Re } \lambda)_{w^{-1}\alpha} \leq 0$$

because $w^{-1}\alpha \in -\Sigma_\lambda^+$ and $\text{Re } \lambda \in (\mathfrak{a}^*)^+$. Thus $\alpha \notin \Sigma_\lambda^+$. This proves (37). Finally, (38) follows immediately from (37). \hfill \square

As intersection of parabolic subgroups of $W$, for each $\lambda \in \mathfrak{a}_c^*$ the group $W_\lambda$ is itself a parabolic subgroup. By definition, this means that there is $I \subset \Pi$ and $w \in W$ so that $W_\lambda = W_{\Theta(I)} = wW_Iw^{-1}$. The proof of the following proposition provides an explicit way of constructing the elements $w$ and $I$ when $\lambda$ is given.

**Proposition 2.2.** Let $W_1, W_2 \subset W$ be two parabolic subgroups. Then $W_1 \cap W_2$ is parabolic.

**Proof.** (see [24, Proposition 3.11]). We can suppose that $W_1 \neq W_2$. Let $\lambda_1$ and $\lambda_2$ be two distinct elements of $\mathfrak{a}^*$ fixed by $W_1$ and $W_2$, respectively. (Recall that every parabolic subgroup of $W$ is the centralizer of some element of $\mathfrak{a}^*$.) Then $W_1 \cap W_2$ fixes the segment $\lambda_1\lambda_2$. Recall also that $\mathfrak{a}_c^* = \sqcup_{w \in W_1 \cap \Pi}(wC_I)$, where

$$C_I = \{ \lambda \in \mathfrak{a}_c^* : \langle \alpha, \lambda \rangle = 0 \text{ for } \alpha \in I, \langle \alpha, \lambda \rangle > 0 \text{ for } \alpha \notin I \} ;$$

see e.g. [19, Section 1.15]. We can therefore find $I \subset \Pi$, $w \in W$ and $\mu_1 \neq \mu_2$ so that $\mu_1\mu_2 \subset \lambda_1\lambda_2 \cap wC_I$. Hence $\mu_1$ and $\mu_2$ admit the same centralizer $wW_Iw^{-1}$. It follows that $wW_Iw^{-1}$ fixes $\lambda_1$ and $\lambda_2$. Hence $wW_Iw^{-1} \subset W_1 \cap W_2$. Conversely, every element in $W_1 \cap W_2$ fixes $\mu_1, \mu_2 \in \lambda_1\lambda_2$. So $W_1 \cap W_2 \subset wW_Iw^{-1}$. \hfill \square

The following lemma describes the possible singularities of the coefficients of the Harish-Chandra expansion of $\varphi_{\lambda_0}(x)$ for an arbitrarily fixed $\lambda = \lambda_0 \in \mathfrak{a}_c^*$. Notice that, by $W$-invariance, we can always suppose that $\text{Re } \lambda_0 \in (\mathfrak{a}^*)^\perp$. Some parts of this lemma must have been considered by previous authors studying Harish-Chandra expansions. As we could not find references for them, we include their proof for the sake of completeness.

**Lemma 2.3.** Let $\lambda_0 \in \mathfrak{a}_c^*$ with $\text{Re } \lambda_0 \in (\mathfrak{a}^*)^\perp$, and keep the above notation. Let $w \in W$ and $\mu \in 2\Lambda \setminus \{0\}$.

(a) The singularities of the function $c(w\lambda)\Gamma_\mu(w\lambda)$ are at most simple poles along the hyperplanes

$$\mathcal{H}_{\alpha,n} \quad \text{with } \alpha \in \Sigma_\lambda^+ \text{ and } n \in \mathbb{Z}.$$ 

The hyperplane $\mathcal{H}_{\alpha,n}$ is a possible singular hyperplane passing through $\lambda_0$ if and only if $\alpha \in \Sigma_{\lambda_0}$ and $n = \langle \lambda_0 \rangle_\alpha$.

(b) The possible singularities of $c(w\lambda)$ at $\lambda = \lambda_0$ are at most simple poles along the hyperplanes

$$\mathcal{H}_{\alpha,0} \quad \text{with } \alpha \in \Sigma_{\lambda_0}^0,$$

$$\mathcal{H}_{\alpha,n} \quad \text{with } \alpha \in \Sigma_{\lambda_0}^0 \cap w(-\Sigma_\lambda^+) \text{ and } n = \langle \lambda_0 \rangle_\alpha \in \mathbb{N}.$$ 

In fact, each hyperplane $\mathcal{H}_{\alpha,0}$ with $\alpha \in \Sigma_{\lambda_0}^0$ is always a simple pole of $c(w\lambda)$ at $\lambda = \lambda_0$. \hfill 13
The possible singularities of $\Gamma_\mu(w\lambda)$ at $\lambda = \lambda_0$ are at most simple poles along the hyperplanes

$$\mathcal{H}_{\alpha,n} \quad \text{with } \alpha \in \Sigma^+_{\lambda_0} \cap w(\Sigma^+_{\lambda_0}) \text{ and } n = (\lambda_0)_{\alpha} \in \mathbb{N}.$$  

(c) Suppose $w \in W_{\text{Re} \lambda_0}$. Then the singularities of $c(w\lambda)$ at $\lambda = \lambda_0$ are precisely simple poles along the hyperplanes

$$\mathcal{H}_{\alpha,0} \quad \text{with } \alpha \in \Sigma^0_{\lambda_0}.$$ 

Those of $\Gamma_\mu(w\lambda)$ at $\lambda = \lambda_0$ are at most first order poles along the hyperplanes

$$\mathcal{H}_{\alpha,n} \quad \text{with } \alpha \in \Sigma^0_{\lambda_0} \text{ and } n = (\lambda_0)_{\alpha} \in \mathbb{N}.$$  

Proof. Because of Theorem 1.2 and Lemma 1.9 the possible singularities of $c(w\lambda)\Gamma_\mu(w\lambda)$ are at most first order poles along hyperplanes of the form

$$\{ \lambda \in a^*_C : w\lambda \in \mathcal{H}_{\alpha,n} \} = w^{-1}\mathcal{H}_{\alpha,n}$$

with $\alpha \in \Sigma^0_\lambda$ and $n \in \mathbb{Z}$. Notice that $w^{-1}\mathcal{H}_{\alpha,n} = \mathcal{H}_{-w^{-1}\alpha,n}$ and that $\mathcal{H}_{-\beta,-m} = \mathcal{H}_{\beta,m}$. The possible singular hyperplanes of $c(w\lambda)\Gamma_\mu(w\lambda)$ are hence of the form given in (a). The last statement in (a) is immediate from the definition of $\mathcal{H}_{\alpha,n}$.

To prove (b), observe that $\mathcal{H}_{\alpha,n}$ is a possible singular hyperplane of $c(w\lambda)$ at $\lambda_0$ if and only if $\mathcal{H}_{w^{-1}\alpha,n} = w^{-1}\mathcal{H}_{\alpha,n}$ is a possible singular hyperplane of $c(\lambda)$ and $\lambda_0 \in \mathcal{H}_{\alpha,n}$. The latter condition is equivalent to $\alpha \in \Sigma^+_{\lambda_0} = \Sigma^+_{\lambda_0} \cup \Sigma^\epsilon_{\lambda_0}$ and $n = (\lambda_0)_{\alpha}$. If $\alpha \in \Sigma^+_{\lambda_0}$, then $n = 0$ and $\mathcal{H}_{w^{-1}\alpha,0}$ is automatically a singular hyperplane of $c(\lambda)$. If $\alpha \in \Sigma^\epsilon_{\lambda_0}$, then $n = (\lambda_0)_{\alpha} \in \mathbb{N}$. In this case $\mathcal{H}_{w^{-1}\alpha,n} = \mathcal{H}_{-w^{-1}\alpha,n}$ is a possible singular hyperplane of $c(\lambda)$ if and only if $-w^{-1}\alpha \in \Sigma^+_{\lambda_0}$. Suppose now $\alpha \in \Sigma^0_{\lambda_0}$ and choose $\beta \in \Sigma^+_{\lambda_0}$ so that $w\alpha \in \{ \pm \beta \}$. Because of the term $\Gamma((w\lambda)\beta)$ at the numerator of the factor $c_\beta(w\lambda)$, the function $c(w\lambda)$ indeed admits a simple pole along $\{ \lambda \in a^*_C : (w\lambda)_{\beta} = 0 \} = \{ \lambda \in a^*_C : \lambda_{w^{-1}\beta} = 0 \} = \mathcal{H}_{\alpha,0}$. 

For $\Gamma_\mu(w\lambda)$, the hyperplane $\mathcal{H}_{\alpha,n}$ is a possible singular hyperplane at $\lambda_0$ if and only if $\mathcal{H}_{w^{-1}\alpha,n} = w^{-1}\mathcal{H}_{\alpha,n}$ is a possible singular hyperplane of $\Gamma_\mu(\lambda)$ and $\lambda_0 \in \mathcal{H}_{\alpha,n}$. This is equivalent to saying that $\mathcal{H}_{w^{-1}\alpha,n}$ is a possible singular hyperplane of $\Gamma_\mu(\lambda)$ and $n = (\lambda_0)_{\alpha} \in \mathbb{N}$, i.e. that $w^{-1}\alpha \in \Sigma^+_{\lambda_0}$, $\alpha \in \Sigma^\epsilon_{\lambda_0}$ and $n = (\lambda_0)_{\alpha}$.

Part (c) is a consequence of (b), (37) and (38). □

To compute the exponential series expansion of the hypergeometric function $\varphi_{\lambda_0}(x)$, we shall need some elementary facts on polynomial differential operators. We collect them in the following lemma whose proof is straightforward.

Lemma 2.4. Let $\lambda_0 \in a^*_C$, and let $I \subset a^*_C \times a^*_C$ be a finite set so that $(\lambda_0 - \nu_2, \nu_1) = 0$ for all $(\nu_1, \nu_2) \in I$. Then the following properties hold.

(a) Define polynomial functions $p_I$ and $\pi_I$ on $a^*_C$ by

$$p_I(\lambda) := \prod_{(\nu_1, \nu_2) \in I} \langle \lambda - \nu_2, \nu_1 \rangle,$$

$$\pi_I(\lambda) := \prod_{(\nu_1, \nu_2) \in I} \langle \lambda, \nu_1 \rangle.$$  

So $p_I(\lambda) = \pi_I(\lambda) + \tilde{p}_I(\lambda)$ with $\deg \tilde{p}_I < \deg p_I = \deg \pi_I = |I|$. Then:

1) $\partial(\pi_I)(p) = 0$ if $p \in S(a_C)$ and $\deg p < |I|$.  

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2) $\partial(p(p_I))_{\lambda=\lambda_0} = 0$ if $p \in S(a_C)$ and $\deg p < |I|$,
3) $\partial(\pi_I(p_I))_{\lambda=\lambda_0} = \partial(\pi_I)(p_I) > 0$.

(b) For every differentiable function $f$ on $a_C^*$ and every $x \in a$ we have

$$
\partial(\pi_I)(f(x)e^{z\lambda(x)})_{\lambda=\lambda_0} = \sum_{j \cup L = I} (\partial(\pi_j)f)(x)_{\lambda=\lambda_0} \pi_L(w^{-1}x)e^{z\lambda_0(x)}
$$

where

$$
\pi_L(x) := \prod_{(\nu_1, \nu_2) \in L} \nu_1(x).
$$

We now construct the polynomials $p(\lambda)$ and $\pi(\lambda)$ appearing on the right-hand side of (41). For $\alpha \in \Sigma_{\lambda_0}$ we set $n_\alpha = (\lambda_0)_\alpha \in \mathbb{N}_0$. Let $w \in W$. Define polynomial functions $\pi_0(\lambda)$, $\pi_1(\lambda)$, $p_{w,+}(\lambda)$, $p_{w,-}(\lambda)$, $p_1(\lambda)$ and $p(\lambda) \in S(a_C)$ by

$$
\pi_0(\lambda) = \prod_{\alpha \in \Sigma_{\lambda_0}} \langle \lambda, \alpha \rangle,
$$

$$
\pi_1(\lambda) = \prod_{\alpha \in \Sigma_{\lambda_0}} \langle \lambda, \alpha \rangle,
$$

$$
\pi(\lambda) = \pi_0(\lambda)\pi_1(\lambda) = \prod_{\alpha \in \Sigma_{\lambda_0}} \langle \lambda, \alpha \rangle,
$$

$$
p_{w,+}(\lambda) = \prod_{\alpha \in \Sigma_{\lambda_0} \cap w(\Sigma_{\lambda_0}^+)} (\langle \lambda, \alpha \rangle - n_\alpha \langle \alpha, \alpha \rangle),
$$

$$
p_{w,-}(\lambda) = \prod_{\alpha \in \Sigma_{\lambda_0} \cap w(-\Sigma_{\lambda_0}^+)} (\langle \lambda, \alpha \rangle - n_\alpha \langle \alpha, \alpha \rangle),
$$

$$
p_1(\lambda) = p_{w,+}(\lambda)p_{w,-}(\lambda),
$$

$$
p(\lambda) = \pi_0(\lambda)p_1(\lambda).
$$

We adopt the convention that empty products are equal to the constant 1. Notice that $p_{w,+}(\lambda)p_{w,-}(\lambda)$ is in fact independent of $w \in W$ and that

$$
p(\lambda) = p_1(\lambda) \quad \text{and} \quad \pi(\lambda) = \pi_1(\lambda) \quad \text{for} \quad I = \{(\alpha, n_{\alpha} \alpha) : \alpha \in \Sigma_{\lambda_0}, n_{\alpha} = (\lambda_0)_\alpha \}.
$$

**Proposition 2.5.** Keep the assumptions of Lemma 2.3.

(a) There is a neighborhood $U$ of $\lambda_0$ with compact closure $\overline{U}$ so that $\overline{U} \cap H_{\alpha,n} \neq \emptyset$ if and only if $\alpha \in \Sigma_{\lambda_0}$ and $n = (\lambda_0)_\alpha$.

(b) For all $w \in W$ and $\mu \in 2\Lambda \setminus \{0\}$, the functions $\pi_0(\lambda)p_{w,-}(\lambda)c(\mu\lambda)$ and $p_{w,+}(\lambda)\Gamma_{\mu}(w\lambda)$ are holomorphic in a neighborhood of $\overline{U}$.

(c) For all $x \in A_C$ where $\varphi_{\lambda_0}(x)$ is defined, we have

$$
c_0\varphi_{\lambda_0}(x) = \partial(\pi)(p(\lambda)\varphi_{\lambda}(x))_{\lambda=\lambda_0}
$$

where $c_0 = \partial(\pi)(p) = \partial(\pi)(\pi) > 0$. 

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(d) Let \( x_0 \in \mathfrak{a}^+ \) be fixed. Then
\[
(52) \quad c_0 \varphi_{\lambda_0}(x) = \sum_{\mu \in 2\Lambda} \sum_{w \in W} \partial(\pi) \left( p(\lambda) c(w\lambda) \Gamma_\mu(w\lambda)e^{(w\lambda - \rho - \mu)(x)} \right)_{\lambda = \lambda_0}
\]
where the series on the right-hand side converges uniformly in \( x \in x_0 + \mathfrak{a}^+ \).

Proof. Parts (a) and (b) are immediate consequences of Lemma 2.3 and the fact that the hyperplanes \( \mathcal{H}_{\beta,n} (\beta \in \Sigma_0^+, n \in \mathbb{Z}) \) form a locally finite family.

To prove (c), notice first that by parts 1) and 2) of Lemma 2.4 (a), we have for \( J \subset I \)
\[
\partial(\pi_J)(p_I)_{\lambda = \lambda_0} = \delta_{J,I} \partial(\pi_I)(\pi_I),
\]
where \( \delta_{J,I} \) is Kronecker’s delta. Hence parts (b) and (c) of the same lemma give for \( I \) as in (50):
\[
\partial(\pi) \left( p(\lambda) \varphi_\lambda(x) \right)_{\lambda = \lambda_0} = \sum_{J \subset I \subset \Lambda} \left( \partial(\pi_J)(p)(\lambda) \right)_{\lambda = \lambda_0} \left( \partial(\pi_I)(\varphi_\lambda(x)) \right)_{\lambda = \lambda_0}
= \partial(\pi)(\pi) \varphi_{\lambda_0}(x).
\]

According to (b) and Lemma 1.4, the series
\[
\sum_{\mu \in 2\Lambda} p_{w,+}(\lambda) \Gamma_\mu(w\lambda)e^{(w\lambda - \rho - \mu)(x)}
\]
converges to \( p_{w,+}(\lambda) \Phi_{w\lambda}(x) \) uniformly in \( U \times (x_0 + \mathfrak{a}^+) \). Moreover, it can be differentiated term-by-term. If \( \lambda \) is regular, then \( \varphi_\lambda(x) = \sum_{w \in W} c(w\lambda) \Phi_{w\lambda}(x) \) for \( x \in \mathfrak{a}^+ \). Multiplying both sides by \( p(\lambda) \), we therefore get for all regular \( \lambda \in U \)
\[
(53) \quad p(\lambda) \varphi_\lambda(x) = \sum_{w \in W} \pi_0(\lambda)p_{w,-}(\lambda)c(w\lambda) \sum_{\mu \in 2\Lambda} p_{w,+}(\lambda) \Gamma_\mu(w\lambda)e^{(w\lambda - \rho - \mu)(x)}.
\]
Since both sides of (53) are holomorphic on \( U \), this equality extends to all of \( U \). Part (d) now follows from (c) and term-by-term differentiation using Lemma 1.4.

A careful computation of (52) is what leads to an expansion of \( \varphi_\lambda(x) \) for all \( \lambda \). We first consider the terms corresponding to \( w \in W_{\lambda_0} \).

Lemma 2.6. Let the notation be as above. Let \( \lambda_0 \in \mathfrak{a}_C^+ \) be such that \( \text{Re} \lambda_0 \in (\mathfrak{a}^+)^\ast \). Define
\[
(54) \quad b_0(\lambda) = \pi_0(\lambda)c(\lambda).
\]
Then the following properties hold for \( w \in W_{\lambda_0} \).

(a) \( \pi_0(w\lambda) = (\det w)\pi_0(\lambda) \) for all \( \lambda \in \mathfrak{a}_C^\ast \).

(b) The point \( \lambda = \lambda_0 \) is neither a zero nor a pole of the function \( b_0(w\lambda) \).

(c) \( b_0(w\lambda_0) = b_0(\lambda_0) \) is a nonzero constant. Moreover, if \( \lambda_0 \in (\mathfrak{a}^+)^\ast \), then \( b_0(\lambda_0) > 0 \).

Proof. Part (a) is classical, as \( W_{\lambda_0} \) is the Weyl group of \( \Sigma_0^0 \sqcup (-\Sigma_0^0) \). All other properties are an immediate consequence of the defining formula (19) of the c-function, the properties of the gamma function, the fact that \( w\lambda_0 = \lambda_0 \) and the assumption \( \text{Re} \lambda_0 \in (\mathfrak{a}^+)^\ast \).
According to (c) in Lemma 2.3, multiplication by \( \pi_0 \) removes all the singularities of \( c(w\lambda) \) at \( \lambda_0 \) for all \( w \in W_{\lambda_0} \). To compute the corresponding terms in (52), we first rewrite them in terms of the function \( b_0 \). Indeed, for \( w \in W_{\lambda_0} \) we have by Lemma 2.6 (a):

\[
(55) \quad \partial_w \left( b_0(w\lambda)p_1(\lambda)e^{(w\lambda-\rho)(x)} \right) |_{\lambda=\lambda_0} = \\
= \left\{ b_0(\lambda_0) \left[ \sum_{J \cup I = J, |J| = |\Sigma_{\lambda_0}^0|} \partial(\pi_J)(p_1) |_{\lambda=\lambda_0} \pi_{wL}(x) \right] + f_{w,\lambda_0}(x) \right\} e^{(\lambda_0-\rho)(x)}
\]

The terms in (52) which correspond to \( w \in W_{\lambda_0} \) and \( \mu = 0 \) are then given by the following lemma. Recall the set \( I \) introduced in (50) and recall that \( \Gamma_0 = 1 \).

**Lemma 2.7.** Let \( w \in W_{\lambda_0} \). Then for \( x \in \mathfrak{a} \) we have

\[
(56) \quad \partial(\pi)(b_0(w\lambda)p_1(\lambda)e^{(w\lambda-\rho)(x)}) |_{\lambda=\lambda_0} = \\
\sum_{J \cup I = J, |J| = |\Sigma_{\lambda_0}^0|} \partial(\pi_J)(p_1) |_{\lambda=\lambda_0} \partial(\pi_L)(b_0(w\lambda)e^{(w\lambda-\rho)(x)}) |_{\lambda=\lambda_0}.
\]

Notice that \( \pi_L(w^{-1}x) = \pi_{wL}(x) \). Applying again Lemma 2.4 (b), we get

\[
\partial(\pi_L)(b_0(w\lambda)e^{(w\lambda-\rho)(x)}) |_{\lambda=\lambda_0} = \\
= \sum_{R \cup S = L} \partial(\pi_R)b_0(w\lambda) |_{\lambda=\lambda_0} \pi_S(w^{-1}x)e^{(\lambda_0-\rho)(x)}
\]

The polynomial \( \pi_{wL}(x) \) has maximal degree (equal to \( |\Sigma_{\lambda_0}^0| \)) if \( |J| = |\Sigma_{\lambda_0}^>^\prime| \). Collecting together these terms proves then the required formula.

To sum the contributions of the terms corresponding to \( \mu = 0 \) and \( w \in W_{\lambda_0} \), we still need two more lemmas.

**Lemma 2.8.** For \( L \subset I \) with \( |L| = |\Sigma_{\lambda_0}^0| \), define

\[
(57) \quad f_L(x) = \sum_{w \in W_{\lambda_0}} (\det w)\pi_{wL}(x), \quad x \in \mathfrak{a}.
\]

Then there is a constant \( c_L \) so that \( f_L(x) = c_L\pi_0(x) \) for all \( x \in \mathfrak{a} \).

**Proof.** The polynomial \( f_L \) is of the same degree as \( \pi_0 \) and it is a \( W_{\lambda_0} \)-skew-symmetric, so divisible by \( \pi_0 \). See e.g. [10, Lemma 10] for a proof of the latter fact. \( \square \)
Lemma 2.9. Let $c_L$ denote the constants introduced in Lemma 2.8 and keep the notation from Lemma 2.7. Set

$$\rho_0 = \sum_{\alpha \in \Sigma^0_{\lambda_0}} \alpha.$$  \hspace{1cm} (58)

Then

$$\sum_{J \cup L = I \atop |J| = |\Sigma^0_{\lambda_0}|} c_L \partial(\pi_J)(p_1)\big|_{\lambda = \lambda_0} = \sum_{J \cup L = I \atop |J| = |\Sigma^0_{\lambda_0}|} c_L \partial(\pi_J)(\pi_1)\big|_{\lambda = \lambda_0} = \frac{1}{\pi_0(\rho_0)} \partial(\pi)(\pi).$$ \hspace{1cm} (59)

Proof. The first equality in (59) follows from Lemma 2.4, (a), part 2). To prove the second equality, notice first that, by Lemmas 2.4 and 2.8, we have

$$\partial(\pi)\left[ \pi_1(\lambda) \left( \sum_{w \in W_{\lambda_0}} (\det w)e^{w\lambda(x)} \right) \right]_{\lambda=0} = \sum_{w \in W_{\lambda_0}} (\det w) \partial(\pi)(\pi_1(\lambda)e^{w\lambda(x)})\big|_{\lambda=0}$$

$$= \sum_{w \in W_{\lambda_0}} (\det w) \sum_{J \cup L = I \atop |J| = |\Sigma^0_{\lambda_0}|} \partial(\pi_J)(\pi_1)\pi_wL(x)$$

$$= \sum_{J \cup L = I \atop |J| = |\Sigma^0_{\lambda_0}|} \partial(\pi_J)(\pi_1)f_L(x)$$

$$= \pi_0(x) \sum_{J \cup L = I \atop |J| = |\Sigma^0_{\lambda_0}|} c_L \partial(\pi_J)(\pi_1).$$

Choose $x = x_{\rho_0}$. Then

$$\sum_{w \in W_{\lambda_0}} (\det w)e^{w\lambda(x_{\rho_0})} = \sum_{w \in W_{\lambda_0}} (\det w)e^{w\rho_0(x_\lambda)} = \prod_{\alpha \in \Sigma^0_{\lambda_0}} \sinh(\alpha, \lambda).$$

The last equality is e.g. Proposition 5.15 (i) in [17, Ch. II, §5]. Since $\pi_0(x_{\rho_0}) = \pi_0(\rho_0) > 0$, we can write

$$\sum_{J \cup L = I \atop |J| = |\Sigma^0_{\lambda_0}|} c_L \partial(\pi_J)(\pi_1) = \frac{1}{\pi_0(\rho_0)} \partial(\pi)\left[ \pi_1(\lambda) \prod_{\alpha \in \Sigma^0_{\lambda_0}} \sinh(\alpha, \lambda) \right]\big|_{\lambda=0}.$$

We claim that

$$\partial(\pi)\left[ \pi_1(\lambda) \prod_{\alpha \in \Sigma^0_{\lambda_0}} \sinh(\alpha, \lambda) \right]\big|_{\lambda=0} = \partial(\pi)(\pi).$$ \hspace{1cm} (60)

Indeed

$$\partial(\pi)\left[ \pi_1(\lambda) \prod_{\alpha \in \Sigma^0_{\lambda_0}} \sinh(\alpha, \lambda) \right]\big|_{\lambda=0} = \sum_{R \cup S = I \atop |R| = |\Sigma^0_{\lambda_0}|} \partial(\pi_R)(\pi_1)\left[ \partial(\pi_S)\left( \prod_{\alpha \in \Sigma^0_{\lambda_0}} \sinh(\alpha, \lambda) \right) \right]\big|_{\lambda=0}.$$
Because of the evaluation at \( \lambda = 0 \), the only nonzero terms inside the last square parenthesis are those of the form

\[
\prod_{j=1}^{d} \partial((\lambda, \beta_j))(\sinh(\gamma_j, \lambda))\big|_{\lambda=0}
\]

where \( d = |\Sigma_0^0| \), and \( \{\beta_1, \ldots, \beta_d\} \) and \( \{\gamma_1, \ldots, \gamma_d\} \) are respectively enumerations of \( S \) and \( \Sigma_0^0 \). The conclusion follows as

\[
\partial((\lambda, \beta))(\sinh(\gamma, \lambda))\big|_{\lambda=0} = (\gamma, \beta) \cosh(\gamma, \lambda)\big|_{\lambda=0} = (\gamma, \beta),
\]

so \( \partial(\pi_S(\prod_{\alpha \in \Sigma_0^0} \sinh(\alpha, \lambda)))\big|_{\lambda=0} = \partial(\pi_S(\pi_0)) \).

**Corollary 2.10.** Keep the above notation. Then

\[
\sum_{w \in W_{\lambda_0}} \partial(\pi(p(\lambda)c(w\lambda)e^{(w\lambda-\rho)(x)})\big|_{\lambda=\lambda_0} = \left(\frac{c_0}{\pi_0(\rho_0)}b_0(\lambda_0)\pi_0(x) + f_{\lambda_0}(x)\right)e^{(\lambda_0-\rho)(x)}
\]

where \( c_0 = \partial(\pi_0)(x) \) is the positive constant of Proposition 2.7 (c), \( b_0(\lambda_0) \neq 0 \) and \( f_{\lambda_0}(x) \) is a polynomial function of \( x \) of degree \( \deg \pi_0 = |\Sigma_0^0| \).

**Proof.** This is an immediate consequence of (55) and of Lemmas 2.6, 2.7, 2.8 and 2.9. \( \square \)

Similar (but less explicit) computations can be performed to evaluate each term of the series (52). Notice that, by Lemma 2.3 (b) and (c), the function defined by

\[
b_w(\lambda) = \begin{cases} 
\pi_0(\lambda)c(w\lambda) & \text{if } w \in W_{Re \lambda_0} \\
\pi_0(\lambda)p_{w,-}(\lambda)c(w\lambda) & \text{if } w \in W \setminus W_{Re \lambda_0}
\end{cases}
\]

is always holomorphic on a neighborhood of \( \mathcal{U} \ni \lambda_0 \). It is nonzero at \( \lambda = \lambda_0 \) for \( w \in W_{Re \lambda_0} \). However, Lemma 2.6 (a) holds only when \( w \in W_{\lambda_0} \). As in the case of \( W_{\lambda_0} \), each term of (52) corresponding to a fixed \( w \in W \) and \( \mu = 0 \) includes a polynomial factor in \( x \). It is of degree \( \deg \pi_0 \) if \( w \in W_{Re \lambda_0} \), and of degree \( \deg \pi_0 + |\Sigma_0^\gamma_+| \) if \( w \in W \setminus W_{Re \lambda_0} \). Estimates for each of the terms of (52) can be obtained from Lemma 1.4. The result of this computation is presented in the following theorem.

**Theorem 2.11.** Keep the assumptions of Proposition 2.5, and let \( x_0 \in a^+ \) be fixed. Then for \( x \in x_0 + a^+ \) we have

\[
(\varphi_0(\lambda_0)) = \left(\frac{c_0}{\pi_0(\rho_0)}b_0(\lambda_0)\pi_0(x) + f_{\lambda_0}(x)\right)e^{(\lambda_0-\rho)(x)}
\]

\[
+ \sum_{w \in (W_{Re \lambda_0} \setminus W_{\lambda_0}) \cup (W \setminus W_{Re \lambda_0})} \left(b_w(\lambda_0)\pi_w,\lambda_0(x) + f_{w,\lambda_0}(x)\right)e^{(w\lambda_0-\rho)(x)}
\]

\[
+ \sum_{\mu \in 2\lambda \setminus \{0\}} \sum_{w \in W} f_{w,\mu,\lambda_0}(x)e^{(w\lambda_0-\rho-\mu)(x)}.
\]

The first term in (63) is as in Corollary 2.10. For \( w \in W \setminus W_{\lambda_0} \), the constant \( b_w(\lambda_0) \) is given by evaluation of (62) at \( \lambda = \lambda_0 \). It is nonzero for \( w \in W_{Re \lambda_0} \setminus W_{\lambda_0} \). The polynomial
\( \pi_{w, \lambda_0}(x) \) is explicitly given by

\[
\pi_{w, \lambda_0}(x) = \begin{cases} 
\sum_{J \cup L = I} \frac{\partial (\pi_J)(\pi_1)}{|J| = |\Sigma^+_{\lambda_0} \cap w(\Sigma^+_{\lambda_0})|} \pi_{wL}(x) & \text{if } w \in W_{\Re \lambda_0} \setminus W_{\lambda_0} \\
\sum_{J \cup L = I} \frac{\partial (\pi_J)(\pi_{w,+})}{|J| = |\Sigma^+_{\lambda_0} \cap w(\Sigma^+_{\lambda_0})|} \pi_{wL}(x) & \text{if } w \in W \setminus W_{\Re \lambda_0}
\end{cases}
\]

where \( \pi_1(\lambda) \) is as in (64) and

\[
\pi_{w,+}(\lambda) = \prod_{\alpha \in \Sigma^+_{\lambda_0} \cap w(\Sigma^+_{\lambda_0})} \langle \lambda, \alpha \rangle.
\]

Moreover, \( f_{w, \lambda_0}(x) \) is a polynomial function of \( x \) with \( \deg f_{w, \lambda_0} < \deg \pi_{w, \lambda_0} \).

For \( \mu \in 2\Lambda \setminus \{0\} \) and \( w \in W \), \( f_{w, \mu, \lambda_0}(x) \) is a polynomial function of \( x \) of degree \( \leq |\Sigma_{\lambda_0}| = \deg p \). The series on the right-hand side of (63) converges uniformly for \( x \in x_0 + \overline{a^+} \).

**Remark 2.12.** The convergence of the series (63) uses the following estimate, which is a consequence of Lemma 1.4. Let \( x_1 \in a^+ \) be fixed. Then there is a constant \( M_{x_1} > 0 \) so that for all \( \mu \in 2\Lambda \setminus \{0\} \) and \( w \in W \) we have

\[
|f_{w, \mu, \lambda_0}(x)| \leq M_{x_1}(1 + |x|)^{|\Sigma_{\lambda_0}|} e^{u(x_1)}, \quad x \in x_1 + \overline{a^+},
\]

This will also be needed in the following section, to determine estimates of the series at infinity in \( a^+ \) away from its walls.

The following corollary gives some relevant special cases of Theorem 2.11

**Corollary 2.13.** Let \( \lambda_0 \in a^+_C \) with \( \Re \lambda_0 \in (a^*)^+ \). If not stated otherwise, we keep assumptions and notation of Theorem 2.11.

(a) Suppose that \( \lambda_0 \) is generic, i.e. \( \lambda_0, \alpha \notin \mathbb{Z} \) for all \( \alpha \in \Sigma_0^+ \). Then for \( x \in x_0 + \overline{a^+} \) we have

\[
\varphi_{\lambda_0}(x) = \sum_{\lambda \in W_{\Re \lambda_0}} c(w_{\lambda_0}) e^{(w_{\lambda_0} - \rho)(x)} + \sum_{\mu \in 2\Lambda \setminus \{0\}} \sum_{w \in W} f_{w, \mu, \lambda_0} e^{(w_{\lambda_0} - \rho - \mu)(x)}
\]

where \( c(w_{\lambda_0}) \neq 0 \) for \( w \in W_{\Re \lambda_0} \) and \( f_{w, \mu, \lambda_0} \in \mathbb{C} \) for \( w \in W \) and \( \mu \in 2\Lambda \setminus \{0\} \).

(b) Suppose \( \langle \lambda_0, \alpha \rangle \neq 0 \) for all \( \alpha \in \Sigma_0^+ \). Then for \( x \in x_0 + \overline{a^+} \) we have

\[
\varphi_{\lambda_0}(x) = \sum_{\lambda \in W_{\Re \lambda_0}} c(w_{\lambda_0}) e^{(w_{\lambda_0} - \rho)(x)} +
\]

\[
+ \frac{1}{c_0} \sum_{w \in W \setminus W_{\Re \lambda_0}} \left( b_{w, \lambda_0}(x) \pi_{w, \lambda_0}(x) + f_{w, \lambda_0}(x) e^{(w_{\lambda_0} - \rho)(x)} \right) +
\]

\[
+ \frac{1}{c_0} \sum_{\mu \in 2\Lambda \setminus \{0\}} \sum_{w \in W} f_{w, \mu, \lambda_0}(x) e^{(w_{\lambda_0} - \rho - \mu)(x)}
\]

with \( c(w_{\lambda_0}) \neq 0 \) for all \( w \in W_{\Re \lambda_0} \).
(c) Suppose that \( \text{Im} \lambda_0 \) belongs to the subspace of \( \mathfrak{a}^* \) supporting the facet containing \( \text{Re} \lambda_0 \), i.e. \( \langle \text{Im} \lambda_0, \alpha \rangle = 0 \) for all \( \alpha \in \Sigma^+_0 \) whenever \( \langle \text{Re} \lambda_0, \alpha \rangle = 0 \). This happens for instance if \( \text{Im} \lambda_0 = 0 \). Then for \( x \in x_0 + \overline{\mathfrak{a}^*} \) we have

\[
(69) \quad c_0 \varphi_{\lambda_0}(x) = \left( \frac{c_0}{\pi_0(\rho_0)} b_0(\lambda_0) \pi_0(x) + f_{\lambda_0}(x) \right) e^{(\lambda_0 - \rho)(x)} + \sum_{w \in W \setminus W_{\text{Re} \lambda_0}} (b_w(\lambda_0) \pi_{w,\lambda_0}(x) + f_{w,\lambda_0}(x)) e^{(w \lambda_0 - \rho)(x)}
+ \sum_{\mu \in 2\Lambda \setminus \{0\}} \sum_{w \in W} f_{w,\mu,\lambda_0}(x) e^{(w \lambda_0 - \rho - \mu)(x)}.
\]

**Proof.** If \( \lambda_0 \) is generic, then \( \Sigma_{\lambda_0} = \emptyset \) and all polynomials in (43) to (49) reduce to the constant polynomial 1. In this case, \( \partial(p) \) is the identity operator and \( c_0 = 1 \). For all \( w \in W \) the functions \( c(w \lambda) \Gamma_\mu(w \lambda) \) are nonsingular at \( \lambda_0 \), and (69) reduces to the (holomorphic extension of the) classical Harish-Chandra expansion we started with.

If \( \langle \lambda_0, \alpha \rangle \neq 0 \) for all \( \alpha \in \Sigma^+_0 \), then \( W_{\lambda_0} = \{ \text{id} \} \) and \( \pi_0 \) is the constant polynomial 1. Hence \( b_0(\lambda_0) = c(\lambda_0) \) and \( b_w(\lambda_0) = c(w \lambda_0) \) for \( w \in W_{\text{Re} \lambda_0} \). Moreover, the polynomials \( f_{\lambda_0} \) and \( f_{w,\lambda_0} \) (with \( w \in W_{\text{Re} \lambda_0} \)) are identically zero. Finally, for \( w \in W_{\text{Re} \lambda_0} \), the polynomials \( \pi_w,\lambda_0 \) are constants and \( \pi_{w,\lambda_0}(x) = \partial(\pi_1(\pi_1) = \partial(p)(p) = c_0 \).

The reduction in (c) follows as \( W_{\text{Re} \lambda_0} \subset W_{\text{Im} \lambda_0} \). So \( W_{\lambda_0} = W_{\text{Re} \lambda_0} \). \qed

3. **Estimates of the hypergeometric functions**

We begin this section by examining the behavior of the hypergeometric functions at infinity on \( \mathfrak{a}^+ \). Our basic tool is the exponential series expansion of \( \varphi_\lambda(x) \) from Theorem 2.11. However, as in the classical case of spherical functions, one cannot work with this series close to the walls of \( \mathfrak{a}^+ \), but on regions of the form \( x_0 + \overline{\mathfrak{a}^*} \) for a fixed \( x_0 \in \mathfrak{a}^+ \). When \( \lambda \in \mathfrak{a}^* \), the remaining region in \( \overline{\mathfrak{a}^+} \) can be handled by means of a subadditivity property proven for \( \varphi_\lambda(x) \) by Schapira in [26].

We keep the notation of the previous section. As in [9] p. 164, define \( \beta : \mathfrak{a} \to \Bbb{R} \) by

\[
(70) \quad \beta(x) = \min_{\alpha \in \Pi} \alpha(x),
\]

where \( \Pi = \{ \alpha_1, \ldots, \alpha_l \} \) denotes as before the set of simple roots in \( \Sigma^+ \). Let \( x_0 \in \mathfrak{a}^+ \) be fixed. Then \( |x| \asymp \beta(x) \) as \( x \to \infty \) in \( x_0 + (\mathfrak{a}^*)^+ \).

Recall that if \( \text{Re} \lambda_0 \in (\mathfrak{a}^*)^+ \) and \( w \in W \), then \( -w \text{Re} \lambda_0 + \text{Re} \lambda_0 = \sum_{j=1}^l r_j \alpha_j \) with \( r_j \geq 0 \). Hence

\[
(71) \quad (\text{Re} \lambda_0 - w \text{Re} \lambda_0)(x) \geq r_w \beta(x), \quad x \in \mathfrak{a}^+,
\]

where \( r_w = \sum_{j=1}^l r_j \geq 0 \). Moreover, \( r_w > 0 \) if \( w \notin W_{\text{Re} \lambda_0} \). We will say that \( x \in \mathfrak{a}^+ \) tends to infinity away from the walls if \( \alpha(x) \to +\infty \) for all \( \alpha \in \Sigma^+_0 \), i.e. if \( \beta(x) \to +\infty \).
Theorem 3.1. Let $\lambda_0 \in a^*$ with $\text{Re} \lambda_0 \in (a^*)^+$, and let $x_0 \in a^+$ be fixed. Then there are constants $C_1 > 0$, $C_2 > 0$ and $b > 0$ (depending on $\lambda_0$ and $x_0$) so that for all $x \in x_0 + \overline{a^+}$:

\[ (72) \quad \frac{\varphi_{\lambda_0}(x)e^{-\text{Re} \lambda_0(x)}}{\pi_0(x)} - \left( \frac{b_0(\lambda_0)}{\pi_0(\rho_0)} e^{i\text{Im} \lambda_0(x)} + \sum_{w \in W_{\text{Re} \lambda_0} \setminus W_{\lambda_0}} \frac{b_w(\lambda_0)\pi_{w,\lambda_0}(x)}{c_0\pi_0(x)} e^{iw\text{Im} \lambda_0(x)} \right) \leq C_1(1 + \beta(x))^{-1} + C_2(1 + \beta(x))|\Sigma_{\lambda_0}^0| e^{-b\beta(x)}. \]

The term $C_1(1 + \beta(x))^{-1}$ on the right-hand side of (72) does not occur if $\Sigma_{\lambda_0}^0 = \emptyset$.

Proof. According to Theorem 2.11, the left-hand-side of (72) is bounded by

\[ \frac{|f_{\lambda_0}(x)|}{c_0\pi_0(x)} + \sum_{w \in W_{\text{Re} \lambda_0} \setminus W_{\lambda_0}} \frac{|f_{w,\lambda_0}(x)|}{c_0\pi_0(x)} + \sum_{w \in W \setminus W_{\text{Re} \lambda_0}} \left| \frac{b_w(\lambda_0)}{c_0\pi_0(x)} \right| |\pi_{w,\lambda_0}(x)| \leq C_f(1 + |x|)^{d_f} \]

for $x \in x_0 + \overline{a^+}$. This leads to the estimate of the right-hand side of (72). By (71) and since $\deg \pi_{w,\lambda_0}(x) = d$, there is a constant $M > 0$ so that for all $w \in W \setminus W_{\text{Re} \lambda_0}$ and $x \in x_0 + \overline{a^+}$ we have

\[ \sum_{w \in W \setminus W_{\text{Re} \lambda_0}} \frac{|b_w(\lambda_0)||\pi_{w,\lambda_0}(x)|}{c_0\pi_0(x)} \leq Me^{-r\beta(x)} \]

with $r = \min_{w \in W \setminus W_{\text{Re} \lambda_0}} t_w > 0$.

For the last part of the estimate we proceed similarly to [9, p. 163]. Apply (66) with $x_1 = x_0/2$. Since $(w\text{Re} \lambda_0 - Re \lambda_0)(x) \leq 0$, we have

\[ \sum_{\mu \in 2\Lambda \setminus \{0\}} \sum_{w \in W} \frac{|f_{w,\mu,\lambda_0}(x)|}{c_0\pi_0(x)} e^{(w\text{Re} \lambda_0 - Re \lambda_0 - \mu)(x)} \leq M'(1 + \beta(x))|\Sigma_{\lambda_0}^{\mu}| d \sum_{\mu \in 2\Lambda \setminus \{0\}} e^{-\mu(nx_0/2)} \leq M'(1 + \beta(x))|\Sigma_{\lambda_0}^{\mu}| \sum_{\mu \in \Lambda \setminus \{0\}} e^{-\mu(2x_0-x)} \]

where $M'$ is a positive constant. Observe that if $x \in x_0 + \overline{a^+}$ then $2x - x_0 \in x_0 + \overline{a^+}$. Recall from Section 1.2. the notation $\ell(\mu) = \sum_{j=1}^{\ell} H_j$ for the level of $\mu = \sum_{j=1}^{\ell} \mu_j \alpha_j \in \Lambda$. For every $m \in \mathbb{N}$ there are at most $m^{\ell-1}$ distinct $\mu \in \Lambda$ with $\ell(\mu) = m$. For $X \in x_0 + \overline{a^+}$ we have then

\[ \sum_{\mu \in \Lambda \setminus \{0\}} e^{-\mu(X)} \leq \sum_{m=1}^{\infty} m^{\ell-1} e^{-m\beta(X)} = e^{-\beta(X)} \sum_{n=0}^{\infty} (n + 1)^{\ell-1} e^{-n\beta(X)}. \]
If \( \beta(X) \geq 1 \), then
\[
\sum_{\mu \in \Lambda \setminus \{0\}} e^{-\mu(X)} \leq e^{-|\beta(X)|} \sum_{n=0}^{\infty} (n+1)^{l-1} e^{-n}.
\]
Since \( \beta(2x-x_0) \geq 1 \) for \( x \in x_0 + \mathbb{a}^+ \) and since \( \beta(2x-x_0) \geq 2|\beta(x) - \max_{j=1,\ldots,l} \alpha_j(x_0)| \), we conclude that
\[
(1 + \beta(x))^{\sum_{\alpha \in \Lambda_0}^\infty} \sum_{\mu \in 2\Lambda \setminus \{0\}} e^{-\mu(x-x_0/2)} \leq C_{x_0} (1 + \beta(x))^{\sum_{\alpha \in \Lambda_0}^\infty} e^{-2|\beta(x)|}.
\]
\[
\square
\]

The next corollary restates Theorem 3.1 in the special case where \( \lambda_0 \in \langle \mathbb{a}^* \rangle^+ \).

**Corollary 3.2.** Let \( \lambda_0 \in \langle \mathbb{a}^* \rangle^+ \), and let \( x_0 \in \mathbb{a}^+ \) be fixed. Then there are constants \( C_1 > 0 \), \( C_2 > 0 \) (depending on \( \lambda_0 \) and \( x_0 \)) so that for all \( x \in x_0 + \mathbb{a}^+ \):
\[
(73) \quad \left| \varphi_{\lambda_0}(x) - \frac{b_0(\lambda_0)}{\pi_0(\rho_0)} \pi_0(x) e^{(\lambda_0-\rho)(x)} \right| \leq
\]
\[
\leq [C_1(1 + \beta(x))^{-1} + C_2(1 + \beta(x))^{\sum_{\alpha \in \Lambda_0}^\infty} e^{-b|\beta(x)|}] \pi_0(x) e^{(\lambda_0-\rho)(x)}.
\]

The term \( C_1(1 + \beta(x))^{-1} \) on the right-hand side of (73) does not occur if \( \sum_{\alpha \in \Lambda_0}^\infty \emptyset \).

Let \( \lambda \in \mathbb{a}^* \). In [26, Lemma 3.4], Schapira proved that for all \( x \in \mathbb{a}^+ \)
\[
(74) \quad \nabla \varphi_{\lambda}(x) = -\frac{1}{|W|} \sum_{\omega \in W} w^{-1}(\rho - \lambda) G_\lambda(w\xi),
\]
where \( G_\lambda \) is the nonsymmetric hypergeometric function. (The gradient is taken with respect to the space variable \( x \in \mathbb{a} \).) Since \( \partial(\xi) F = \langle \nabla F, \xi \rangle \) and because of (24), one obtains for all \( \xi \in \mathbb{a} \)
\[
\partial(\xi) \left( e^{K_\xi \frac{\xi \cdot \lambda}{\xi \cdot \xi}} \varphi_\lambda(\cdot) \right) \leq 0
\]
where \( K_\xi = \max_{w \in W} w(\rho - \lambda)(w\xi) \). This in turn yields the following subadditivity property, which is implicit in [26].

**Lemma 3.3.** Let \( \lambda \in \mathbb{a}^* \). Then for all \( x, x_1 \in \mathbb{a} \) we have
\[
(75) \quad \varphi_\lambda(x + x_1) e^{-\max_{w \in W}(\lambda - \rho)(w\xi_1)} \leq \varphi_\lambda(x) \leq \varphi_\lambda(x + x_1) e^{\max_{w \in W}(\lambda - \rho)(w\xi_1)}.
\]
In particular, if \( \lambda \in \langle \mathbb{a}^* \rangle^+ \) and \( x_1 \in \mathbb{a}^+ \), then
\[
(76) \quad \varphi_\lambda(x + x_1) e^{-(\lambda + \rho)(x_1)} \leq \varphi_\lambda(x) \leq \varphi_\lambda(x + x_1) e^{(\lambda + \rho)(x_1)}
\]
for all \( x \in \mathbb{a} \).

Together with Corollary 3.2, the above lemma yields the following global estimates of \( \varphi_\lambda \), stated without proof in [26, Remark 3.1].

**Theorem 3.4.** Let \( \lambda_0 \in \langle \mathbb{a}^* \rangle^+ \). Then for all \( x \in \mathbb{a}^+ \) we have
\[
(77) \quad \varphi_{\lambda_0}(x) \leq \prod_{\alpha \in \sum_{\alpha \in \Lambda_0}^\infty (1 + \alpha(x))} e^{(\lambda_0 - \rho)(x)}.
\]
Proof. We proceed as in [26, Theorem 3.1] or in [1], for the case \( \lambda_0 = 0 \). Choose \( x_1 \in x_0 + (a^*)^+ \) so that the expression inside the square bracket at the right-end side of (73) is smaller than \( \frac{b_0(\lambda_0)}{2\pi(\rho_0)} \) for \( x \in x_1 + (a^*)^+ \). Applying then (76), we obtain the required result from (73). \( \square \)

4. Bounded hypergeometric functions

Let \( C(\rho) \) denote the convex hull of the points \( w \rho (w \in W) \). The main result of this section is Theorem 4.2 below, which extends the theorem of Helgason and Johnson [18] to the hypergeometric functions of Heckman and Opdam. We first point out the following lemma.

Lemma 4.1. We have \( \varphi_\rho \equiv 1 \).

Proof. From (23) we have \( G_{-\rho} \equiv 1 \) and so (24) implies that \( \varphi_{-\rho} \equiv 1 \). If \( w_0 \) is the longest element in \( W \) then \( w_0 \rho = -\rho \) and the \( W \)-invariance of \( \varphi_\lambda \) in the \( \lambda \)-variable implies that \( \varphi_{w \rho} \equiv 1 \) for all \( w \in W \). \( \square \)

Theorem 4.2. The hypergeometric function \( \varphi_\lambda \) is bounded if and only if \( \lambda \) belongs to the tube \( C(\rho) + i a^* \). Moreover, \( |\varphi_\lambda(x)| \leq 1 \) for all \( \lambda \in C(\rho) + i a^* \) and \( x \in a \).

Proof. First we show that \( |\varphi_\lambda(x)| \leq 1 \) if \( \lambda \in C(\rho) + i a^* \). For this, we apply maximum modulus principle to the holomorphic function \( \lambda \mapsto \varphi_\lambda(x) \) (for a fixed \( x \)).

Let \( R > 0 \) be arbitrary but fixed and let \( B_R \) be the closed ball of radius \( R \) centered at the origin in \( a^* \). Applying the maximum modulus principle along with \( |\varphi_\lambda(x)| \leq \varphi_{\Re \lambda}(x) \) in the domain \( C(\rho) + iB_R \) implies that the maximum of \( |\varphi_\lambda(x)| \) is obtained when \( \lambda \) belongs to the boundary of \( C(\rho) \subset a^* \). Let \( \mu_1 \) and \( \mu_2 \) be two distinct points on the boundary of \( C(\rho) \). Let \( \overline{\mu_1,\mu_2} \) be the line segment in \( a^* \) joining \( \mu_1 \) and \( \mu_2 \), and let \( L_{\mu_1,\mu_2}^C \) be the complex line passing through \( \mu_1 \) and \( \mu_2 \), that is,

\[
L_{\mu_1,\mu_2}^C = \{ z \mu_1 + (1-z) \mu_2 : z \in \mathbb{C} \}.
\]

Consider the (closed) domain

\[
P_\rho = (C(\rho) + iB_R) \cap \{ \lambda \in L_{\mu_1,\mu_2}^C : \Re \lambda \in \overline{\mu_1,\mu_2} \}.
\]

Now, a similar argument to the above shows that the maximum of \( \lambda \mapsto |\varphi_\lambda(x)| \) in \( P_\rho \) is attained at either \( \mu_1 \) or \( \mu_2 \). It immediately follows that the maximum of \( |\varphi_\lambda(x)| \) in \( C(\rho) + i a^* \) is attained at the extreme points of \( C(\rho) \) which is just the set \( \{ wp : w \in W \} \). This completes the proof as \( \varphi_{w \rho}(x) \equiv 1 \).

We now prove that for \( \lambda_0 \) such that \( \Re \lambda_0 \notin C(\rho) \) the function \( \varphi_{\lambda_0} \) is not bounded. By \( W \)-invariance in the spectral parameter, we can suppose that \( \Re \lambda_0 \notin (a^*)^+ \). We first recall that

\[
C(\rho) \cap (a^*)^+ = (\rho - (a^*)^+) \cap (a^*)^+(\text{see [17 Ch. IV, Lemma 8.3 (i) \].})
\]

If \( \Re \lambda_0 \notin (a^*)^+ \setminus C(\rho) \), we can therefore find \( x_1 \in a^+ \) so that \( \Re (\lambda_0 - \rho)(x_1) > 0 \). Set \( d = |\Sigma_{\lambda_0}^0| \in \mathbb{N} \). If \( \varphi_{\lambda_0} \) is bounded, we have

\[
\lim_{t \to +\infty} \varphi_{\lambda_0}(tx_1)e^{-t(\Re \lambda_0 - \rho)(x_1)}t^{-d} = 0.
\]

24
Recall that the polynomials $\pi_0$ and $\pi_{w,\lambda_0}$ appearing in Theorem 3.1 are of degree $d$. Observe also that $\pi_0(x_1) \neq 0$ as $x_1 \in \mathfrak{a}^+$. Replacing $x$ by $tx_1$ in (72), we obtain as $t \to +\infty$

$$
\left| \frac{\varphi_{\lambda_0}(tx_1)e^{-t(\text{Re} \lambda_0 - \rho)(x_1)}}{t^d\pi_0(x_1)} - \left( \frac{b_0(\lambda_0)}{\pi_0(\rho_0)} e^{it\text{Im} \lambda_0(x_1)} \right)
\right|
+ \sum_{w \in W_{\text{Re} \lambda_0} \setminus W_{\lambda_0}} \frac{b_w(\lambda_0)\pi_{w,\lambda_0}(x_1)}{c_0\pi_0(x_1)} e^{itw\text{Im} \lambda_0(x_1)}
= o(t).
$$

It follows that

$$(80) \quad \lim_{t \to +\infty} \left( \frac{b_0(\lambda_0)}{\pi_0(\rho_0)} e^{it\text{Im} \lambda_0(x_1)} + \sum_{w \in W_{\text{Re} \lambda_0} \setminus W_{\lambda_0}} \frac{b_w(\lambda_0)\pi_{w,\lambda_0}(x_1)}{c_0\pi_0(x_1)} e^{itw\text{Im} \lambda_0(x_1)} \right) = 0.
$$

We now proceed as in [9 p. 147]. If $u_1, \ldots, u_p$ are distinct complex numbers and $c_1, \ldots, c_p$ are complex constants, then

$$(81) \quad \lim_{t \to +\infty} \left| \sum_{j=1}^{p} c_j e^{iu_jt} \right|^2 \geq \frac{1}{T} \int_0^T \left| \sum_{j=1}^{p} c_j e^{iu_jt} \right|^2 dt = \sum_{j=1}^{p} |c_j|^2.
$$

So $\lim_{t \to +\infty} \left| \sum_{j=1}^{p} c_j e^{iu_jt} \right|^2 = 0$ implies $c_j = 0$ for $j = 1, \ldots, p$. Observe that if $w \in W_{\text{Re} \lambda_0} \setminus W_{\lambda_0}$ then $w \notin W_{\text{Im} \lambda_0}$. Since $x_1 \in \mathfrak{a}^+$, this implies that $w \text{Im} \lambda_0(x_1) \neq \text{Im} \lambda_0(x_1)$ for all $w \in W_{\text{Re} \lambda_0} \setminus W_{\lambda_0}$. Consequently, (80) contradicts that $b_0(\lambda_0) \neq 0$. Thus $\varphi_{\lambda_0}$ cannot be bounded.

**Remark 4.3.** Some results towards Theorem 4.2 have been obtained in [25, Theorem 5.4 and Corollary 5.6]. More precisely, the proof in [25] yields:

(a) $\varphi_\lambda$ is bounded if $\lambda$ belongs to the interior of the tube $C(\rho) + i\mathfrak{a}^*$. (In [25], this is an application of Schapira’s sharp estimates proven in Theorem 3.4 above.)

(b) Suppose $\lambda \notin C(\rho) + i\mathfrak{a}^*$. If, moreover, either $\lambda \in \mathfrak{a}^*$ or $\lambda \in \mathfrak{a}_c^* \setminus \mathfrak{a}^*$ is generic, then $\varphi_\lambda$ is not bounded. (This result is a combination of Schapira’s estimates on $\mathfrak{a}^*$ and the classical Harish-Chandra series for generic $\lambda \in \mathfrak{a}_c^*$.)

Notice that (b) yields that $\varphi_\lambda$ is not bounded if $\lambda \notin C(\rho) + i\mathfrak{a}^*$ in the rank-one case.

Notice also that the original proof by Helgason and Johnson is based on a detailed study of Harish-Chandra’s integral formula for the spherical functions and of the boundary components of the symmetric space $G/K$. These objects are missing in the general theory of hypergeometric functions associated with root systems. So our proof provides an alternative proof of the characterization of the bounded spherical functions as well.

5. **$L^p$-Fourier Analysis**

In this section we present some results towards a development of the $L^p$-harmonic analysis for the hypergeometric Fourier transform $\mathcal{F}$. We begin by considering the holomorphic properties of $\mathcal{F}$ on $L^1(\mathfrak{a}, d\mu)^W$. This is a simple application of Theorem 4.2. A Riemann-Lebesgue lemma is also obtained. We then study the properties of $\mathcal{F}$ on $L^p(\mathfrak{a}, d\mu)^W$ with $1 < p < 2$. We establish Hausdorff-Young inequalities and as a consequence, using an argument from [7 pages 249–250], we obtain injectivity and an inversion formula for $\mathcal{F}$ on $L^p(\mathfrak{a}, d\mu)^W$. Furthermore, we introduce the $L^p$-Schwartz space for $0 < p \leq 2$ and prove an isomorphism theorem. Many of these results are quite similar to those of the geometric case.
the Hausdorff-Young inequalities are proven as in [5]; the $L^p$-Schwartz space isomorphism is obtained using Anker’s method [2]. A crucial ingredient is given by the following estimates, due to Schapira (see [26, Theorem 3.4 and Remark 3.2]): let $p \in S(\mathfrak{a}_C^*)$ and $q \in S(\mathfrak{a}_C)$ be polynomials of degrees respectively $M$ and $N$. Then there is a constant $C > 0$ such that

\begin{equation}
|\partial_\lambda(p)\partial_q(q)\varphi_\lambda(x)| \leq C(1 + |x|)^M(1 + |\lambda|)^N\varphi_0(x)e^{\max_{\omega \in W} \Re w\lambda(x)}
\end{equation}

for all $\lambda \in \mathfrak{a}_C^*$ and $x \in \mathfrak{a}$. In (82), we have written $\partial_y$ to indicate that the differential operator acts on the variable $y$. These estimates are obtained in [26] as a consequence of similar estimates for the functions $G_\lambda$.

For $0 < p \leq 2$, set $\epsilon_p = \frac{2}{p} - 1$. Let $C(\epsilon_p \rho)$ be the convex hull in $\mathfrak{a}^*$ of the set $\{\epsilon_p w \rho : w \in W\}$, and let $\mathfrak{a}_p^* = C(\epsilon_p \rho) + i\mathfrak{a}^*$. In particular, for $p = 1$, we have that $\mathfrak{a}_1^* = C(\rho) + i\mathfrak{a}^*$ is precisely the set of parameters $\lambda$ for which $\varphi_\lambda$ is bounded.

**Corollary 5.1.** Let $f \in L^1(\mathfrak{a}, d\mu)^W$. Then the following properties hold.

(a) The hypergeometric Fourier transform $\hat{f}(\lambda)$ is well defined for all $\lambda \in \mathfrak{a}_1^*$, and

\begin{equation}
|\hat{f}(\lambda)| \leq \|f\|_1, \quad \lambda \in \mathfrak{a}_1^*.
\end{equation}

(b) The function $\hat{f}$ is continuous on $\mathfrak{a}_1^*$ and holomorphic in its interior.

(c) (Riemann-Lebesgue lemma) We have

\[ \lim_{\lambda \in \mathfrak{a}_1^*, \Im \lambda \to \infty} |\hat{f}(\lambda)| = 0. \]

**Proof.** The first two properties are immediate consequences of Theorem 4.2, the fact that $\varphi_\lambda$ is holomorphic in $\lambda$, and Morera’s theorem. The Riemann-Lebesgue lemma follows by approximating an arbitrary function in $L^1(\mathfrak{a}, d\mu)^W$ by $W$-invariant compactly supported smooth functions and the Paley-Wiener theorem. □

Next we discuss some properties of $F$ on $L^p(\mathfrak{a}, d\mu)^W$.

**Lemma 5.2.** Let $f \in L^p(\mathfrak{a}, d\mu)^W$. Then the following properties hold.

(a) The hypergeometric Fourier transform $\hat{f}(\lambda)$ is well defined for all $\lambda$ in the interior of $\mathfrak{a}_p^*$. It defines a holomorphic function in the interior of $\mathfrak{a}_p^*$.

(b) (Hausdorff-Young: real case) Let $p, q$ be so that $1 < p < 2$ and $1/p + 1/q = 1$. Then there is a constant $C_p > 0$ so that $\|\hat{f}\|_q := \left( \int_{\mathfrak{a}_p^*} |\hat{f}(\lambda)|^q |e(\lambda)|^{-2} d\lambda \right)^{1/q} \leq C_p\|f\|_p$.

**Proof.** The first part is an immediate application of the estimates (82). The proof of the second part can be obtained by following the methods used in [5] Lemma 8. More precisely, it is an application to the Riesz-Thorin interpolation theorem to the operator $F$, which is of type $(2, 2)$ by the Plancherel theorem and of type $(1, \infty)$ by (83). □

The Hausdorff-Young inequality above can be extended as in [5]. As the proofs are very similar we only give a sketch. Recall that $\Sigma_0$ is the set of elements in $\Sigma$ which are not integral multiples of other elements in $\Sigma$. For $\alpha \in \Sigma_0$ define, $a(\alpha) = m_\alpha + m_{2\alpha}$.

**Lemma 5.3.** Let $f \in L^p(\mathfrak{a}, d\mu)^W, 1 < p < 2$ and $\eta$ be in the interior of $C(\epsilon_p \rho)$. Then the following properties hold.
Next we define an admissible family of analytic operators as follows: Again, by the Hausdorff-Young inequality along with Paley-Wiener estimates for $L^p$-classes.

By the Plancherel theorem, we have

\[ \left( \int_{ia^*} |\hat{f}(\lambda + \eta)|^q |c(\lambda)|^{-2} d\lambda \right)^{1/q} \leq C_{p,q} \|f\|_p. \]

Proof. To prove (a) we follow ideas from [5]. First, we note that, as in [5, Lemma 5] we have

\[ (84) \quad |c(\lambda)|^{-2} \propto \Pi_{\alpha \in \Sigma_0} |\langle \lambda, \alpha \rangle|^2(1 + |\langle \lambda, \alpha \rangle|)^{a(\alpha)}^{-2}. \]

Next we define an admissible family of analytic operators as follows:

For $z \in \mathbb{C}$ such that $0 \leq \Re z \leq \frac{|\beta|}{p}$, let $Y_z(f)$ be the function defined on the measure space $(ia^*, d\nu)$ where $d\nu(\lambda) = \Pi_{\alpha \in \Sigma_0} (1 + |\langle \lambda, \alpha \rangle|)^{a(\alpha)} d\lambda$ by

\[ (85) \quad Y_z(f)(\lambda) = \hat{f}( \eta|\eta| + \lambda ) \Pi_{\alpha \in \Sigma_0} (1 + |\langle \lambda, \alpha \rangle|)^{-1} (\eta|\eta| + \lambda, \alpha). \]

Proceeding as in [5], we see that $Y_z$ is of type $(2,2)$ when $\Re z = 0$ and type $(1,\infty)$ when $\Re z = \frac{|\beta|}{p}$ with admissible bounds. Analytic interpolation then proves (a). The next two results are established exactly as in [5]. We omit the details.

Now we turn to the question of injectivity and an inversion formula for $\mathcal{F}$.

**Theorem 5.4.** Let $f \in L^p(a, d\mu)^W, 1 \leq p \leq 2$. Then $\hat{f} \equiv 0$ implies that $f \equiv 0$ almost everywhere. Moreover, if $f \in L^p(a, d\mu)^W$ and $f \in L^1(a^*, |c(\lambda)|^{-2} d\lambda)$ then

\[ f(x) = \int_{ia^*} \hat{f}(\lambda) \varphi_{-\lambda}(x) |c(\lambda)|^{-2} d\lambda \quad \text{almost everywhere}. \]

Proof. The first part of the proof follows the argument used in [7, Theorem 3.2]. Fix $g \in C_c^\infty(a)^W$. Consider the linear functionals $T_g$ and $\hat{T}_g$ on $L^p(a, d\mu)^W$ defined by

\[ T_g(h) = \int_a h(x) \overline{g(x)} \ d\mu(x), \]

\[ \hat{T}_g(h) = \int_{ia^*} \hat{h}(\lambda) \overline{\hat{g}(\lambda)} |c(\lambda)|^{-2} d\lambda. \]

By the Plancherel theorem, $T_g = \hat{T}_g$ on $L^1(a, d\mu)^W \cap L^2(a, d\mu)^W$, which is a dense subspace of $L^p(a, d\mu)^W$ for $1 \leq p \leq 2$. By Hölder’s inequality we have $|T_g(h)| \leq \|h\|_p \|g\|_q$ if $\frac{1}{p} + \frac{1}{q} = 1$.

Again, by the Hausdorff-Young inequality along with Paley-Wiener estimates for $\hat{g}(\lambda)$, we have

\[ |\hat{T}_g(h)| \leq \|\hat{g}\|_p \|\hat{h}\|_q \leq C_p \|\hat{g}\|_p \|h\|_p. \]
Hence, both \( T_0 \) and \( \hat{T}_g \) are continuous linear functionals on \( L^p(\mathfrak{a},d\mu)^W \) which agree on a dense subspace. It follows that they agree on all of \( L^p(\mathfrak{a},d\mu)^W \). Now, if \( f \in L^p(\mathfrak{a},d\mu)^W \) and \( \hat{f} \equiv 0 \), then \( \hat{T}_g(f) = T_g(f) = 0 \). So we have,

\[
\int_{\mathfrak{a}} f(x) \overline{g(x)} \, d\mu(x) = 0
\]

for all \( g \in C^\infty_c(\mathfrak{a})^W \), which implies that \( f \) vanishes almost everywhere.

To prove the inversion formula, notice that if \( \hat{f} \in L^1(\mathfrak{a}^*, |c(\lambda)|^{-2}d\lambda)^W \), then by the Fubini’s theorem and \( \varphi_\lambda(x) = \varphi_{-\lambda}(x) \) for \( \lambda \in i\mathfrak{a}^* \), we have

\[
\hat{T}_g(f) = \int_{\mathfrak{a}} g(x) \left( \int_{i\mathfrak{a}^*} \hat{f}(\lambda) \varphi_{-\lambda}(x) |c(\lambda)|^{-2} d\lambda \right) \, d\mu(x).
\]

Since \( T_g(f) = \hat{T}_g(f) \) and \( g \) is arbitrary we get the result. \( \square \)

We record the following equality obtained in the proof of Theorem 5.4.

**Corollary 5.5.** Let \( 1 \leq p \leq 2 \). Suppose that \( f \in L^p(\mathfrak{a},d\mu)^W \) and \( g \in C^\infty_c(\mathfrak{a})^W \). Then

\[
\int_{\mathfrak{a}} f(x) \overline{g(x)} \, dx = \int_{i\mathfrak{a}^*} \hat{f}(\lambda) \overline{\hat{g}(\lambda)} |c(\lambda)|^{-2} d\lambda.
\]

For \( 0 < p \leq 2 \), we define the \( L^p \)-Schwartz space \( C^p(\mathfrak{a})^W \) to be the set of all \( C^\infty W \)-invariant functions \( f \) on \( \mathfrak{a} \) such that for each \( N \in \mathbb{N}_0 \) and \( q \in S(\mathfrak{a}_C) \),

\[
\sup_{x \in \mathfrak{a}} (1 + |x|)^N |\varphi_0(x)|^{-n} |\partial(q)f(x)| < \infty.
\]

(86)

It is easy to check that \( C^\infty_c(\mathfrak{a})^W \subset C^p(\mathfrak{a})^W \subset L^p(\mathfrak{a},d\mu)^W \). Hence \( C^p(\mathfrak{a})^W \) is dense in \( L^p(\mathfrak{a},d\mu)^W \). As in the geometric case, it can be shown that \( C^p(\mathfrak{a})^W \) is a Fréchet space with respect to the seminorms defined by the left-hand side of (86). It can also be shown that \( C^{p_1}(\mathfrak{a})^W \subset C^{p_2}(\mathfrak{a})^W \) when \( p_1 \leq p_2 \) and that the inclusion map is continuous.

Let \( \mathcal{S}(\mathfrak{a}_p^*)^W \) be the set of all \( W \)-invariant functions \( F : \mathfrak{a}_p^* \to \mathbb{C} \) which are holomorphic in the interior of \( \mathfrak{a}_p^* \), continuous on \( \mathfrak{a}_p^* \), and satisfy for all \( r \in \mathbb{N}_0 \) and \( s \in S(\mathfrak{a}_C) \)

\[
\sup_{\lambda \in \mathfrak{a}_p^*} (1 + |\lambda|^r) |\partial(s)F(\lambda)| < \infty.
\]

(87)

We note that when \( p = 2 \), this space reduces to the usual Schwartz space of functions on \( i\mathfrak{a}^* \). It is easy to check that \( \mathcal{S}(\mathfrak{a}_p^*)^W \) is a Fréchet algebra under pointwise multiplication and with the topology induced by the seminorms defined by the left-hand side of (87). Moreover, using the Euclidean Fourier transform, one can prove that \( PW(\mathfrak{a}_C^*)^W \) is a dense subalgebra of \( \mathcal{S}(\mathfrak{a}_p^*)^W \).

Recall that the Euclidean Fourier transform of a sufficiently regular function \( f : \mathfrak{a} \to \mathbb{C} \) is defined by

\[
\tilde{f}(\lambda) = \int_{\mathfrak{a}} f(x)e^{\lambda(x)} \, dx, \quad \lambda \in \mathfrak{a}_C^*.
\]

It is an isomorphism between \( C^\infty_c(\mathfrak{a})^W \) and \( PW(\mathfrak{a}_C^*)^W \). It follows that there is an isomorphism \( \mathcal{A} : C^\infty_c(\mathfrak{a})^W \to C^\infty_c(\mathfrak{a})^W \) such that \( \tilde{\mathcal{A}} f(\lambda) = \tilde{f}(\lambda) \) for all \( \lambda \in \mathfrak{a}_C^* \). In the geometric case, \( \mathcal{A} \) is nothing but the Abel transform.
We now consider the Schwartz space isomorphism theorem. This theorem was proved by Schapira [26, Theorem 4.1] for the case $p = 2$ as an adaptation of Anker’s method for the geometric case (see [2]). Anker’s proof in fact extends to the case of the hypergeometric Fourier transform on $\mathcal{C}^p(a)$ with $0 < p \leq 2$. We outline the main steps of the proof for the reader’s convenience.

**Theorem 5.6.** For $0 < p \leq 2$, the hypergeometric Fourier transform is a topological isomorphism between $\mathcal{C}^p(a)^W$ and $\mathcal{S}(a^*_p)^W$.

**Proof.** Using (82) and the estimates for $\varphi_0$ (which are a special case of Theorem 3.4), one can check, as in the geometric case (see e.g. [9, Theorem 7.8.6]), that the hypergeometric Fourier transform $\mathcal{F}$ maps $\mathcal{C}^p(a)^W$ into $\mathcal{S}(a^*_p)^W$ and is continuous. Since $\mathcal{C}^p(a)^W \subset L^2(a, d\mu)^W$, it follows from the Plancherel theorem that $\mathcal{F}$ is injective. We now show that $\mathcal{F}$ maps $\mathcal{C}^p(a)^W$ into $\mathcal{S}(a^*_p)^W$ is surjective and that its inverse map is continuous. For this, using the Paley-Wiener theorem, it is sufficient to prove that given a seminorm $\beta$ of $\mathcal{C}^p(a)^W$, there exists a seminorm $\eta$ of $\mathcal{S}(a^*_p)^W$ and a constant $C > 0$ such that

$$
\beta(f) \leq C \eta(\hat{f})
$$

for all $f \in C_c^{\infty}(a)$. For $r \in \mathbb{N}_0$ and $s \in \mathcal{S}(a_c)$, let

$$
\eta_{r,s}^{(p)}(F) = \sup_{\lambda \in a^*_p} (1 + |\lambda|)^r |\partial(s) F(\lambda)|, \quad F \in \mathcal{S}(a^*_p)^W
$$

and let

$$
\beta(f) = \sup_{x \in a} (1 + |x|)^N \varphi_0(x)^{-\frac{n}{2}} |\partial(q) f(x)|, \quad f \in \mathcal{C}^p(a)^W.
$$

For each positive integer $r$, let $\Gamma_r = \{ x \in a : \rho(x^+) \leq r \}$ be the polar set of $r^{-1}\rho$ (see section 1.5 for the notation). Then each $\Gamma_r$ is convex, compact and $W$-invariant. Moreover, $a$ is the disjoint union of $\Gamma_1$ and $\Gamma_r \setminus \Gamma_{r-1}$, $r \geq 2$.

Let $f \in C_c^{\infty}(a)^W$. Then, using the inversion formula and the estimates (82), one checks that

$$
\sup_{x \in \Gamma_1} (1 + |x|)^N \varphi_0(x)^{-\frac{n}{2}} |\partial(q) f(x)| \leq C_1 \eta_{N_1}^{(2)}(\hat{f})
$$

for some $N_1 \in \mathbb{N}$. Here we have used that $\Gamma_1$ is compact.

Let $u_r \in C_c^{\infty}(a)^W$ be equal to 1 on $\Gamma_{r-1}$, equal to 0 on $a \setminus \Gamma_r$ and such that $u_r$ and all of its derivatives are bounded, uniformly on $r \in \mathbb{N}$. Since $\mathcal{A}$ is an isomorphism, there exists $f_r \in C_c^{\infty}(a)^W$ such that $(1 - u_r) \mathcal{A} f = \mathcal{A} f_r$. From this it follows that $f$ and $f_r$ may differ only inside $\Gamma_r$.

Using again the inversion formula on $f_r$, one observes that

$$
|\partial(q) f_r(x)| \leq C_2 \varphi_0(x) \eta_{N_1}^{(2)}(\hat{f}_r).
$$

We now estimate $\eta_{N_1}^{(2)}(\hat{f}_r)$. Using the relation $\widehat{\mathcal{A}f_r}(\lambda) = \hat{f}_r(\lambda)$ and Euclidean Fourier analysis, we prove that

$$
\eta_{N_1}^{(2)}(\hat{f}_r) \leq C_3 \sum_{k=0}^{N_1} \sup_{x \in a \setminus \Gamma_{r-1}} (1 + |x|)^{l+1} |\nabla_a^k \mathcal{A} f(x)|.
$$

Here $l = \dim a$ and $\nabla_a$ is the gradient operator on $a$. 

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Next we estimate \( \sup_{x \in \Gamma_{r+1} \setminus \Gamma_r} (1 + |x|)^N \varphi_0(x)^{-\frac{2}{p}} |\partial(q) f(x)| \). Note that \( f = f \) on \( \Gamma_{r+1} \setminus \Gamma_r \). Then, using the previous steps and the estimate of \( \varphi_0 \), we get
\[
(89) \quad \sup_{x \in \Gamma_{r+1} \setminus \Gamma_r} (1 + |x|)^N \varphi_0(x)^{-\frac{2}{p}} |\partial(q) f(x)| \leq C_4 r^N N_2 e^{\rho(x)} \eta_{N_1,1}(f_r)
\]
for some \( N_2 \in \mathbb{N} \). By (88), the right-hand side of is dominated by
\[
C_5 \sum_{k=0}^{N_1} \sup_{x \in a^+ \setminus \Gamma_{r-1}} (1 + |x|)^{N_2 + l + 1} e^{\rho(x)} |\nabla^k A f(x)|.
\]
In turn, this is dominated by
\[
C_6 \sum_{m=0}^{N_2 + l + 1} \int_{a^+} (1 + |\lambda|)^{N_1} |\nabla^m \hat{f}(\lambda + \epsilon \rho)| \, d\lambda.
\]
This inequality follows from the fact that for any polynomials \( p \) and \( q \)
\[
p(x) e^{\rho(x)} \partial(q) g(x) = c \int_{a^+} \partial(p) \{ q(\lambda - \epsilon \rho) \, h(-\lambda + \epsilon \rho) \} e^{-\lambda(x)} \, d\lambda
\]
where \( g(x) = \int_{a^+} h(\lambda) e^{-\lambda(x)} \, d\lambda \). Thus
\[
\sup_{x \in \Gamma_{r+1} \setminus \Gamma_r} (1 + |x|)^N \varphi_0(x)^{-\frac{2}{p}} |\partial(q) f(x)| \leq C_7 \sum_{m=0}^{N_2 + l + 1} \sup_{\lambda \in a^+} (1 + |\lambda|)^{N_1 + l + 1} |\nabla^m \hat{f}(\lambda)|.
\]
This completes the proof.

We conclude this section with a remark on the convolution structure of \( L^1(\mathfrak{a}, d\mu)^W \). Suppose that the triple \( (\mathfrak{a}, \Sigma, m) \) is geometric and corresponds to the Riemannian symmetric space of the noncompact type \( G/K \). The space \( L^1(G/K)^K \) of \( K \)-invariant \( L^1 \) functions on \( G/K \) is a commutative Banach algebra with respect to the \( L^1 \)-norm and a natural positive convolution structure. This structure is transferred to \( L^1(\mathfrak{a}, d\mu)^W \) by restriction to \( \mathfrak{a} \equiv \exp \mathfrak{a} \). The maximal ideals of \( L^1(G/K)^K \) are all of the form
\[
M_\lambda = \{ f \in L^1(G/K)^K : \int_G f(g) \varphi_\lambda(g) \, dg = 0 \}
\]
where \( \varphi_\lambda \) is a bounded spherical function. It follows that, in the geometric case, the bounded hypergeometric functions \( \varphi_\lambda \) parametrize the characters of the commutative algebra \( L^1(\mathfrak{a}, d\mu)^W \). For arbitrary triples, to find a positive convolution structure is an important open problem. When the rank is one, this is completely settled by the results of Flensted-Jensen and Koornwinder [7]. When the rank is greater than one, the only result available is due to Rösler [25] where the root system is of the type \( BC \) and multiplicity function ranges in three continuous one-parameter families. In both cases, it has been proven that the bounded hypergeometric functions indeed give the characters of \( L^1(\mathfrak{a}, d\mu)^W \).

References

[1] J.-P. Anker: La forme exacte de l’estimation fondamentale de Harish-Chandra, C. R. Acad. Sci. Paris, Sér. I Math. 305 (9) (1987) 371–374.
[2] J.-P. Anker: The spherical Fourier transform of rapidly decreasing functions. A simple proof of a characterization due to Harish-Chandra, Helgason, Trombi, and Varadarajan, J. Funct. Anal. 96 (2) (1991), 331–349.
[3] I. Cherednik: A unification of Knizhnik-Zamolodchikov and Dunkl operators via affine Hecke algebras, Invent. Math. 106 (2) (1991), 411-431.
[4] P. Delorme: Transformation de Fourier hypergéométrique, J. Funct. Anal. 168 (1999), 239–312.
[5] M. Eguchi, K. Kunnahara: An $L^p$ Fourier analysis on symmetric spaces, J. Funct. Anal. (2) 47 (1982), 230–246.
[6] A. Erdélyi, et al.: Higher transcendental functions, volume 1, McGraw-Hill, New York, 1953.
[7] M. Flensted-Jensen, T. H. Koornwinder: The convolution structure for Jacobi function expansions, Ark. Mat. 11 (1973), 242–262.
[8] R. Gangolli: On the Plancherel formula and the Paley-Wiener theorem for spherical functions on semisimple Lie groups, Ann. of Math. 93 (1) (1971), 150–165.
[9] R. Gangolli, V. S. Varadarajan: Harmonic analysis of spherical functions on real reductive groups, Springer-Verlag, Berlin, 1988.
[10] Harish-Chandra: Differential operators on a semisimple Lie algebra, Amer. J. Math. 79 (1957), 87–120.
[11] ———: Spherical functions on a semisimple Lie group I, Amer. J. Math. 80 (1958), 241–310.
[12] G. J. Heckman: Root systems and hypergeometric functions II, Compositio Math. 64 (3) (1987), 353–373.
[13] ———: Dunkl operators, In: Séminaire Bourbaki, Vol. 1996/97, Astérisque 245 (1997), Exp. No. 828, 4, pp. 223–246.
[14] G. J. Heckman, E. M. Opdam: Root systems and hypergeometric functions I, Compositio Math. 64 (3) (1987), 329–352.
[15] G. J. Heckman, H. Schlichtkrull: Harmonic analysis and special functions on symmetric spaces, Academic Press, 1994.
[16] S. Helgason: An analogue of the Paley-Wiener theorem for the Fourier transform on certain symmetric spaces, Math. Ann. 165 (1966), 297–308.
[17] ———: Groups and geometric analysis. Integral geometry, invariant differential operators, and spherical functions, Mathematical Surveys and Monographs, 83. AMS, Providence, RI, 2000.
[18] S. Helgason, K. Johnson: The bounded spherical functions on symmetric spaces, Adv. Math. 3 (1969), 586–593.
[19] J. E. Humphreys: Reflection groups and Coxeter groups, Cambridge University Press, 1990.
[20] T. H. Koornwinder: Jacobi functions and analysis on noncompact semisimple Lie groups. In: R. A. Askey, T. H. Koornwinder, W. Schempp, Special functions: group theoretical aspects and applications, Reidel, Dordrecht, 1984, pp. 1–85.
[21] E. M. Opdam: Root systems and hypergeometric functions, IV. Compositio Math. 67 (2) (1988), 191–209.
[22] ———: Harmonic analysis for certain representations of graded Hecke algebras, Acta Math. 175 (1995), 75–12.
[23] ———: Lecture notes on Dunkl operators for real and complex reflection groups, Mathematical Society of Japan, Tokyo, 2000.
[24] D. Qi: On irreducible, infinite, affine Coxeter groups, Ph.D. Dissertation, The Ohio State University, 2007.
[25] M. Rösler: Positive convolution structure for a class of Heckman-Opdam hypergeometric functions of type $BC$, J. Funct. Anal. 258 (8) (2010), 2779–2800.
[26] B. Schapira: Contributions to the hypergeometric function theory of Heckman and Opdam: sharp estimates, Schwartz space, heat kernel, Geom. Funct. Anal. 18 (1) (2008), 222–250.
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