# Uniform oscillatory behavior of spherical functions of $G L_{n} / U_{n}$ at the identity and a central limit theorem 

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# Uniform oscillatory behavior of spherical functions of $G L_{n} / U_{n}$ at the identity and a central limit theorem 

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#### Abstract

Let $\mathbb{F}=\mathbb{R}$ or $\mathbb{C}$ and $n \in \mathbb{N}$. Let $\left(S_{k}\right)_{k \geq 0}$ be a time-homogeneous random walk on $G L_{n}(\mathbb{F})$ associated with an $U_{n}(\mathbb{F})$-biinvariant measure $\nu \in M^{1}\left(G L_{n}(\mathbb{F})\right)$. We derive a central limit theorem for the ordered singular spectrum $\sigma_{\operatorname{sing}}\left(S_{k}\right)$ with a normal distribution as limit with explicit analytic formulas for the drift vector and the covariance matrix. The main ingredient for the proof will be a oscillatory result for the spherical functions $\varphi_{i \rho+\lambda}$ of $\left(G L_{n}(\mathbb{F}), U_{n}(\mathbb{F})\right)$. More precisely, we present a necessarily unique mapping $m_{1}: G \rightarrow \mathbb{R}^{n}$ such that for some constant $C$ and all $g \in G, \lambda \in \mathbb{R}^{n}$,


$$
\left|\varphi_{i \rho+\lambda}(g)-e^{i \lambda \cdot m_{1}(g)}\right| \leq C\|\lambda\|^{2}
$$

KEYWORDS: Biinvariant random walks on $G L(n, \mathbb{R})$ and $G L(n, \mathbb{C})$, asymptotics of spherical functions, central limit theorem for the singular spectrum, random walks on the positive definite matrices, dispersion.

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## 1 Introduction

Let $\mathbb{F}=\mathbb{R}$ or $\mathbb{C}, n \geq 2$ an integer, and $G:=G L(n, \mathbb{F})$ the general linear group with maximal compact subgroup $K:=U_{n}(\mathbb{F})$. Consider i.i.d. $G$-valued random variables $\left(X_{k}\right)_{k \geq 1}$ with the common $K$-biinvariant distribution $\nu \in M^{1}(G)$ and the associated $G$-valued random walk $\left(S_{k}:=X_{1} \cdot X_{2} \cdots X_{k}\right)_{k \geq 0}$ with the convention that $S_{0}$ is the identity $I_{n}$. Moreover, let

$$
\sigma_{\text {sing }}(g) \in\left\{x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}: x_{1} \geq x_{2} \geq \cdots \geq x_{n}>0\right\}
$$

denote the singular (or Lyapunov) spectrum of $g \in G$ where the singular values of $g$, i.e., square roots of the eigenvalues of $g g^{*}$, are ordered by size. Consider the mapping $\ln \sigma_{\text {sing }}$ from $G$ onto the Weyl chamber

$$
W_{n}:=\left\{x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}: x_{1} \geq x_{2} \geq \cdots \geq x_{n}\right\}
$$

with the logarithm $\ln \left(x_{1}, \ldots, x_{n}\right):=\left(\ln x_{1}, \ldots, \ln x_{n}\right)$. We show that under a natural moment condition, the $\mathbb{R}^{n}$-valued random variables

$$
\begin{equation*}
\frac{1}{\sqrt{k}}\left(2 \cdot \ln \sigma_{\text {sing }}\left(S_{k}\right)-k \cdot m_{\mathbf{1}}(\nu)\right) \tag{1.1}
\end{equation*}
$$

tend for $k \rightarrow \infty$ to some $n$-dimensional normal distribution $N\left(0, \Sigma^{2}(\nu)\right)$ where the drift vector $m_{1}(\nu)$ and the covariance matrix $\Sigma^{2}(\nu)$ are given explicitely depending on $\nu$.

This central limit theorem (CLT) can be also seen as follows: By polar decomposition of $g \in G$, the symmetric space $G / K$ can be identified with the cone $P_{n}(\mathbb{F})$ of positive definite symmetric or hermitian $n \times n$ matrices via

$$
g K \mapsto I(g):=g g^{*} \in P_{n}(\mathbb{F}) \quad(g \in G)
$$

where $G$ acts on $P_{n}(\mathbb{F})$ via $a \mapsto g a g^{*}$. In this way, the double coset space $G / / K$ can be identified with the Weyl chamber $W_{n}$ via

$$
K g K \mapsto \ln \sigma_{\text {sing }}(g)=\frac{1}{2} \ln \sigma\left(g g^{*}\right)
$$

where $\sigma$ denotes the spectrum, i.e., the ordered eigenvalues, of a positive definite matrix. Therefore, the CLT above may be regarded as a CLT for the spectrum of $K$-invariant random walks on $P_{n}(\mathbb{F})$. Such CLTs have a long history. CLTs where $\nu$ is renormalized first into some measure $\nu_{k} \in M^{1}(G)$, and then the convergence of the convolution powers $\nu_{k}^{k}$ is studied, can be found e.g. in [KTS], [Tu], [FH], [Te1], [Te2], [Ri], [G1], and [G2]. In this case, so-called dispersions of $\nu$ appear as parameters of the limits, where these dispersions are defined in terms of derivatives of the spherical functions of $(G, K)$. These dispersions will also appear in our CLT in order to describe $m_{1}(\nu)$ and $\Sigma^{2}(\nu)$. Our CLT is in principle well-known; see Theorem 1 of [Vi], as well as the CLTs of Le Page [L] and the monograph [BL]. However, our approach, which directly leads to analytic formulas for drift and covariance, seems to be new for $n>2$. For $n=2$, our CLT can be splitted into two one-dimensional parts, namely a classical part for the sum $\ln \operatorname{det} S_{k}=\sum_{l=1}^{k} \ln \operatorname{det} X_{l}$ of i.i.d. random variables, and a CLT for $\left(S L_{2}(\mathbb{F}), S U(2, \mathbb{F})\right)$. The associated spherical functions are the Jacobi functions $\varphi_{\lambda}^{(0,0)}(t)$ and $\varphi_{\lambda}^{(1 / 2,1 / 2)}(t)$ depending on $\mathbb{F}$ (see $[\mathrm{K}]$ for details), and the CLT above for $\left(S L_{2}(\mathbb{F}), S U(2, \mathbb{F})\right)$ appears as a special case of a CLT of Zeuner [Z] for certain Sturm-Liouville hypergroups on $[0, \infty[$. The proof of Zeuner depends on some uniform estimate for the oscillatory behavior of the associated multiplicative functions, i.e., the Jacobi functions here. This idea was later on transfered to certain random walks on the nonnegative integers associated with orthogonal polynomials in [V1]. Moreover, the result of Zeuner [Z] was recently slightly improved for Jacobi functions in [V2] by using the well-known Harish-Chandra integral representation of the Jacobi functions from $[\mathrm{K}]$. We here adopt this approach and use the Harish-Chandra integral representation of the spherical function of $(G, K)$ to derive a uniform estimate for their oscillatory behavior. The CLT above then follows easily.

Let us describe this uniform oscillatory result. Recapitulate that a $K$-biinvariant continuous function $\varphi \in C(G)$ on $G$ is called spherical iff

$$
\varphi\left(g_{1}\right) \varphi\left(g_{2}\right)=\int_{K} \varphi\left(g_{1} k g_{2}\right) d k
$$

for all $g_{1}, g_{2} \in G$ where $d k$ is the normalized Haar measure on $K$. It is well-known (see [H1] or $[\mathrm{Te} 2])$ that all spherical functions of $(G, K)$ are given by the Harish-Chandra integral

$$
\begin{equation*}
\varphi_{i \rho+\lambda}(g)=\int_{K} \Delta_{1}^{i \lambda_{1}-i \lambda_{2}}\left(k^{*} g g^{*} k\right) \cdots \Delta_{n-1}^{i \lambda_{n-1}-i \lambda_{n}}\left(k^{*} g g^{*} k\right) \Delta_{n}^{i \lambda_{n}}\left(k^{*} g g^{*} k\right) d k \tag{1.2}
\end{equation*}
$$

where the $\Delta_{r}$ are the principal minors of order $r, \lambda \in \mathbb{C}^{n}$, and where $\rho=\left(\rho_{1}, \ldots, \rho_{n}\right)$ is the half sum of roots with $\rho_{l}=\frac{d}{2}(n+1-2 l)$ with the dimension $d=1,2$ of $\mathbb{F}$ over $\mathbb{R}$. Notice that by (1.2), $\varphi_{i \rho} \equiv 1$, and that for $\lambda \in \mathbb{R}^{n}$ and $g \in G,\left|\varphi_{i \rho+\lambda}(g)\right| \leq 1$.

We now follow the usual approach to the dispersion for $(G, K)$ (see $[\mathrm{FH}],[\mathrm{Te} 1],[\mathrm{Te} 2],[\mathrm{Ri}]$, [G1], [G2]) and to so-called moment functions on hypergroups in [Z], [V1], and Section 7.2.2 of $[\mathrm{BH}]$ : For multiindices $l=\left(l_{1}, \ldots, l_{n}\right) \in \mathbb{N}_{0}^{n}$ we define the so called moment functions

$$
\begin{align*}
m_{l}(g) & :=\left.\frac{\partial^{|l|}}{\partial \lambda^{l}} \varphi_{i \rho-i \lambda}(g)\right|_{\lambda=0}:=\left.\frac{\partial^{|l|}}{\left(\partial \lambda_{1}\right)^{l_{1}} \cdots\left(\partial \lambda_{n}\right)^{l_{n}}} \varphi_{i \rho-i \lambda}(g)\right|_{\lambda=0} \\
& =\int_{K}\left(\ln \Delta_{1}\left(k^{*} g g^{*} k\right)\right)^{l_{1}} \cdot\left(\ln \left(\frac{\Delta_{2}\left(k^{*} g g^{*} k\right)}{\Delta_{1}\left(k^{*} g g^{*} k\right)}\right)\right)^{l_{2}} \cdots\left(\ln \left(\frac{\Delta_{n}\left(k^{*} g g^{*} k\right)}{\Delta_{n-1}\left(k^{*} g g^{*} k\right)}\right)\right)^{l_{n}} d k \tag{1.3}
\end{align*}
$$

of order $|l|:=l_{1}+\cdots+l_{n}$ for $g \in G$. Clearly, the last equality follows immediately from (1.2) by interchanging integration and derivatives. Using the $n$ moment functions $m_{l}$ of first order $|l|=1$, we form the vector-valued moment function

$$
\begin{equation*}
m_{\mathbf{1}}(g):=\left(m_{(1,0, \ldots, 0)}(g), \ldots, m_{(0, \ldots, 0,1)}(g)\right) \tag{1.4}
\end{equation*}
$$

of first order. Moreover, we use the usual scalar product $x \cdot y:=\sum_{l=1}^{n} x_{l} y_{l}$ on $\mathbb{R}^{n}$. We can now formulate the following oscillatory result; it will be proved in Section 2.
1.1 Theorem. There exists a constant $C=C(n)$ such that for all $g \in G$ and and $\lambda \in \mathbb{R}^{n}$,

$$
\left|\varphi_{i \rho+\lambda}(g)-e^{i \lambda \cdot m_{1}(g)}\right| \leq C\|\lambda\|^{2}
$$

The function $m_{1}$ is obviously determined uniquely by the property of the theorem.
We return to the CLT. Similar to collecting the moment functions of first order in the vector $m_{1}$, we group the moment functions of second order by

$$
\begin{align*}
m_{\mathbf{2}}(g) & :=\left(\begin{array}{ccc}
m_{1,1}(g) & \cdots & m_{1, n}(g) \\
\vdots & & \vdots \\
m_{n, 1}(g) & \cdots & m_{n, n}(g)
\end{array}\right)  \tag{1.5}\\
: & =\left(\begin{array}{cccc}
m_{(2,0, \ldots, 0)}(g) & m_{(1,1,0, \ldots, 0)}(g) & \cdots & m_{(1,0, \ldots, 0,1)}(g) \\
m_{(1,1,0, \ldots, 0)}(g) & m_{(0,2,0, \ldots, 0)}(g) & \cdots & m_{(0,1,0, \ldots, 0,1)}(g) \\
\vdots & \vdots & & \vdots \\
m_{(1,0, \ldots, 0,1)}(g) & m_{(0,1,0, \ldots, 0,1)}(g) & \cdots & m_{(0, \ldots, 0,2)}(g)
\end{array}\right)
\end{align*}
$$

for $g \in G$. We show in Section 3 as an easy consequence of (1.3) that the $n \times n$ matrices $m_{\mathbf{2}}(g)-m_{\mathbf{1}}(g)^{t} \cdot m_{\mathbf{1}}(g)$ are positive semidefinite.

Now consider $\nu \in M^{1}(G)$ such that the moment functions $m_{j, j} \geq 0(j=1, \ldots, n)$ are $\nu$-integrable. We then say that $\nu$ admits finite second moments. In this case, (1.3) and the Cauchy-Schwarz inequality yield that all moments of order one and two are $\nu$-integrable, and we form the modified expectation vector and covariance matrix

$$
m_{\mathbf{1}}(\nu):=\int_{G} m_{\mathbf{1}}(g) d \nu \in \mathbb{R}^{n}, \quad \Sigma^{2}(\nu):=\int_{G} m_{\mathbf{2}}(g) d \nu-m_{\mathbf{1}}(\nu)^{t} \cdot m_{\mathbf{1}}(\nu)
$$

of $\nu$. The precise statement of our CLT is now as follows:
1.2 Theorem. If $\nu \in M^{1}(G)$ is $K$-biinvariant and admits finite second moments, then

$$
\begin{equation*}
\frac{1}{\sqrt{k}}\left(2 \cdot \ln \sigma_{\text {sing }}\left(S_{k}\right)-k \cdot m_{\mathbf{1}}(\nu)\right) \longrightarrow N\left(0, \Sigma^{2}(\nu)\right) \tag{1.6}
\end{equation*}
$$

for $k \rightarrow \infty$ in distribution.

This paper is organized as follows: Section 2 is devoted exclusively to the proof of Theorem 1.1. In Section 3 we then shall present the proof of Theorem 1.2. There we also give a precise condition on $\nu$ when $\Sigma^{2}(\nu)$ is positive definite.

We finally remark that the results of our paper can be transfered to the Grassmann manifolds $\left(S O_{0}(p, q) /(S O(p) \times S O(q))\right.$ and $(S U(p, q) / S(U(p) \times U(q))$. In this case, the spherical functions are certain Heckman-Opdam hypergeometric functions of type BC, for which a Harish-Chandra integral representation analog to (1.2) is available; see [Sa] and [RV]. We plan to carry out this in near future.

## 2 Proof of the oscillatory behavior of spherical functions

This section is devoted to the proof of Theorem 1.1 which depends on several facts which may be more or less well-known. As we could not find suitable published references, we include proofs for sake of completeness. We start with a result about the principal minors $\Delta_{r}$ :
2.1 Lemma. Let $1 \leq r \leq n$ be integers, $\mathbb{F}=\mathbb{R}$ or $\mathbb{C}$, and $u \in U_{n}(\mathbb{F})$. Then

$$
\Delta_{r}\left(u^{*} \cdot \operatorname{diag}\left(a_{1}, \ldots, a_{n}\right) \cdot u\right)=\sum_{1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n} c_{i_{1}, \ldots, i_{r}} a_{i_{1}} \cdot a_{i_{2}} \cdots a_{i_{r}}
$$

for all $a_{i_{1}}, a_{i_{2}}, \ldots, a_{i_{r}} \in \mathbb{R}$ with coefficients $c_{i_{1}, \ldots, i_{r}}=c_{i_{1}, \ldots, i_{r}}(u)$ satisfying $c_{i_{1}, \ldots, i_{r}} \geq 0$ for $1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n$ and $\sum_{1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n} c_{i_{1}, \ldots, i_{r}}=1$.

Proof. Clearly, $h_{r}\left(a_{1}, \ldots, a_{n}\right):=\Delta_{r}\left(u^{*} \cdot \operatorname{diag}\left(a_{1}, \ldots, a_{n}\right) \cdot u\right)$ is a homogeneous polynomial of degree $r$, i.e.,

$$
h_{r}\left(a_{1}, \ldots, a_{n}\right)=\sum_{1 \leq i_{1} \leq i_{2} \leq \ldots \leq i_{r} \leq n} c_{i_{1}, \ldots, i_{r}} a_{i_{1}} \cdot a_{i_{2}} \cdots a_{i_{r}} .
$$

We first check that $c_{i_{1}, \ldots, i_{r}} \neq 0$ is possible only for $1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n$. For this consider $i_{1}, \ldots, i_{r}$ with $\left|\left\{i_{1}, \ldots, i_{r}\right\}\right|=: q<r$. By changing the numbering of the variables $a_{1}, \ldots, a_{n}$ (and of rows and columns of $u$ in an appropriate way), we may assume that $\left\{i_{1}, \ldots, i_{r}\right\}=\{1, \ldots, q\}$. In this case, $u^{*} \cdot \operatorname{diag}\left(a_{1}, \ldots, a_{q}, 0, \ldots, 0\right) \cdot u$ has rank at most $q<r$. Thus

$$
0=h_{r}\left(a_{1}, \ldots, a_{q}, 0, \ldots, 0\right)=\sum_{1 \leq i_{1} \leq i_{2} \leq \ldots \leq i_{r} \leq q} c_{i_{1}, \ldots, i_{r}} a_{i_{1}} \cdot a_{i_{2}} \cdots a_{i_{r}}
$$

for all $a_{1}, \ldots, a_{q}$. This yields $c_{i_{1}, \ldots, i_{r}}=0$ for $1 \leq i_{1} \leq i_{2} \leq \ldots \leq i_{r} \leq q$ and proves that

$$
h_{r}\left(a_{1}, \ldots, a_{n}\right)=\sum_{1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n} c_{i_{1}, \ldots, i_{r}} a_{i_{1}} \cdot a_{i_{2}} \cdots a_{i_{r}} .
$$

For the nonnegativity we again may restrict our attention to $c_{1, \ldots, r}$. In this case,

$$
0 \leq\left(\begin{array}{cc}
I_{r} & 0 \\
0 & 0
\end{array}\right) \leq I_{n} \quad \text { and thus } \quad 0 \leq u^{*}\left(\begin{array}{cc}
I_{r} & 0 \\
0 & 0
\end{array}\right) u \leq I_{n}
$$

w.r.t. the usual ordering of positive semidefinite matrices. As this inequality holds also for the upper left $r \times r$ block, we obtain

$$
c_{1, \ldots, r}=h_{r}(1, \ldots, 1,0, \ldots, 0)=\Delta_{r}\left(u^{*}\left(\begin{array}{cc}
I_{r} & 0 \\
0 & 0
\end{array}\right) u\right) \geq 0 .
$$

Finally, as

$$
\sum_{1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n} c_{i_{1}, \ldots, i_{r}}=h_{r}(1, \ldots, 1)=1
$$

the proof is complete.
Let us keep the notation of Lemma 2.1. We now compare $h_{r}\left(a_{1}, \ldots, a_{n}\right)$ with the homogeneous polynomial

$$
\begin{equation*}
C_{r}\left(a_{1}, \ldots, a_{n}\right):=\frac{1}{\binom{n}{r}} \sum_{1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n} a_{i_{1}} a_{i_{2}} \cdots a_{i_{r}}>0 \quad(r=1, \ldots, n) \tag{2.1}
\end{equation*}
$$

2.2 Lemma. For all $a_{1}, \ldots, a_{n}>0$,

$$
0<\frac{C_{r}\left(a_{1}, \ldots, a_{n}\right)}{h_{r}\left(a_{1}, \ldots, a_{n}\right)} \leq \frac{1}{\binom{n}{r}} \sum_{1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n} c_{i_{1}, \ldots, i_{r}}(u)^{-1}
$$

where, depending on $u$, on both sides the value $\infty$ is possible.
Proof. Positivity is clear by Lemma 2.1. Moreover,

$$
\begin{aligned}
C_{r}\left(a_{1}, \ldots, a_{n}\right) & =\frac{1}{\binom{n}{r}} \sum_{1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n} a_{i_{1}} a_{i_{2}} \cdots a_{i_{r}} \\
& \leq \frac{\max _{1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n} c_{i_{1}, \ldots, i_{r}}^{-1}}{\binom{n}{r}} \sum_{1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n} c_{i_{1}, \ldots, i_{r}} a_{i_{1}} a_{i_{2}} \cdots a_{i_{r}}
\end{aligned}
$$

which immediately leads to the claim.
We also need the following observation from linear algebra:
2.3 Lemma. Let $u \in U_{n}(\mathbb{C})$ have the block structure $u=\left(\begin{array}{cc}u_{1} & * \\ * & u_{2}\end{array}\right)$ with quadratic blocks $u_{1} \in M_{r}(\mathbb{C})$ and $u_{2} \in M_{n-r}(\mathbb{C})$ with $1 \leq r \leq n$. Then $\left|\operatorname{det} u_{1}\right|=\left|\operatorname{det} u_{2}\right|$.

Proof. W.l.o.g. we may assume $2 r \leq n$. By the $K A K$-decomposition of $U_{n}(\mathbb{C})$ with $K=$ $U_{r}(\mathbb{C}) \times U_{n-r}(\mathbb{C})$ (see e.g. Theorem VII.8.6 of [H2]), we write $u$ as

$$
u=\left(\begin{array}{cc}
a_{1} & 0 \\
0 & b_{1}
\end{array}\right) \cdot\left(\begin{array}{ccc}
c & s & 0 \\
-s & c & 0 \\
0 & 0 & I_{q-2 r}
\end{array}\right) \cdot\left(\begin{array}{cc}
a_{2} & 0 \\
0 & b_{2}
\end{array}\right)
$$

with $a_{1}, a_{2} \in U_{r}(\mathbb{C}), b_{1}, b_{2} \in U_{n-r}(\mathbb{C})$ and with $c=\operatorname{diag}\left(\cos \varphi_{1}, \ldots, \cos \varphi_{r}\right)$ and $s=$ $\operatorname{diag}\left(\sin \varphi_{1}, \ldots, \sin \varphi_{r}\right)$ for suitable $\varphi_{1}, \ldots, \varphi_{r} \in \mathbb{R}$. Therefore,

$$
u_{1}=a_{1} c a_{2} \quad \text { and } \quad u_{2}=b_{1}\left(\begin{array}{cc}
c & 0 \\
0 & I_{q-2 r}
\end{array}\right) b_{2}
$$

which immediately implies the claim.
We shall also need the following elementary observation:
2.4 Lemma. Let $\varepsilon \in] 0,1], M \geq 1$ and $m \in \mathbb{N}$. Then there exists a constant $C=$ $C(\varepsilon, M, m)>0$ such that for all $z \in] 0, M]$,

$$
|\ln (z)|^{m} \leq C\left(1+z^{-\varepsilon}\right)
$$

Proof. Elementary calculus yields $\left|x^{\varepsilon} \cdot \ln x\right| \leq 1 /(e \varepsilon)$ for $\left.\left.x \in\right] 0,1\right]$ and the Euler number $e=2,71 \ldots$. This leads to the estimate for $z \in] 0,1]$. The estimate is trivial for $z \in] 1, M]$.

Proof of Theorem 1.1: As the spherical functions and the moment functions on $G$ are constant on the double cosets w.r.t. $K$ by definition, and as each double coset has a representative $g$ such that $g g^{*}=\operatorname{diag}\left(a_{1}, \ldots, a_{n}\right)$ is diagonal with $a_{1} \geq \ldots \geq a_{n}>0$, it suffices to consider these group elements $g \in G$. We thus fix $\lambda \in \mathbb{R}^{n}$ and $a_{1} \geq \ldots \geq a_{n}>0$ and put $a:=\operatorname{diag}\left(a_{1}, \ldots, a_{n}\right)$. According to (1.2), (1.3) and (1.4) we have to estimate

$$
\begin{align*}
R & :=R(\lambda, a):=\left|\varphi_{i \rho+\lambda}(g)-e^{i \lambda \cdot m_{\mathbf{1}}(g)}\right|  \tag{2.2}\\
= & \mid \int_{K} \exp \left(i \sum_{r=1}^{n}\left(\lambda_{r}-\lambda_{r+1}\right) \cdot \ln \Delta_{r}\left(k^{*} a k\right)\right) d k \\
\quad & \quad-\exp \left(i \int_{K} \sum_{r=1}^{n}\left(\lambda_{r}-\lambda_{r+1}\right) \cdot \ln \Delta_{r}\left(k^{*} a k\right) d k\right) \mid
\end{align*}
$$

with the convention $\lambda_{n+1}:=0$. For $r=1, \ldots, n$, we now use the polynomial $C_{r}$ from Eq. (2.1) and write the logarithms of the principal minors in (2.2) as

$$
\begin{equation*}
\ln \Delta_{r}\left(k^{*} a k\right)=\ln C_{r}\left(a_{1}, \ldots, a_{r}\right)+\ln \left(H_{r}(k, a)\right) \quad \text { with } \quad H_{r}(k, a):=\frac{\Delta_{r}\left(k^{*} a k\right)}{C_{r}\left(a_{1}, \ldots, a_{n}\right)} \tag{2.3}
\end{equation*}
$$

With this notation and with $\left|e^{i x}\right|=1$ for $x \in \mathbb{R}$, we rewrite (2.2) as

$$
\begin{align*}
R=\mid \int_{K} & \exp \left(i \sum_{r=1}^{n}\left(\lambda_{r}-\lambda_{r+1}\right) \cdot \ln \left(H_{r}(k, a)\right)\right) d k \\
& \quad-\exp \left(i \int_{K} \sum_{r=1}^{n}\left(\lambda_{r}-\lambda_{r+1}\right) \cdot \ln \left(H_{r}(k, a)\right) d k\right) \mid \tag{2.4}
\end{align*}
$$

We now use the power series for both exponential functions and observe that the terms of order 0 and 1 are equal in the difference above. Hence,

$$
R \leq R_{1}+R_{2}
$$

for

$$
R_{1}:=\int_{K}\left|\exp \left(i \sum_{r=1}^{n}\left(\lambda_{r}-\lambda_{r+1}\right) \cdot \ln \left(H_{r}(k, a)\right)\right)-\left(1+i \sum_{r=1}^{n}\left(\lambda_{r}-\lambda_{r+1}\right) \cdot \ln \left(H_{r}(k, a)\right)\right)\right| d k
$$

and
$R_{2}:=\left|\exp \left(i \int_{K} \sum_{r=1}^{n}\left(\lambda_{r}-\lambda_{r+1}\right) \cdot \ln \left(H_{r}(k, a)\right) d k\right)-1-i \int_{K} \sum_{r=1}^{n}\left(\lambda_{r}-\lambda_{r+1}\right) \cdot \ln \left(H_{r}(k, a)\right) d k\right|$.

Using the well-known elementary estimates $|\cos x-1| \leq x^{2} / 2$ and $|\sin x-x| \leq x^{2} / 2$ for $x \in \mathbb{R}$, we obtain $\left|e^{i x}-(1+i x)\right| \leq x^{2}$ for $x \in \mathbb{R}$. Therefore, defining

$$
A_{m}:=\int_{K}\left|\sum_{r=1}^{n}\left(\lambda_{r}-\lambda_{r+1}\right) \cdot \ln \left(H_{r}(k, a)\right)\right|^{m} d k \quad(m=1,2)
$$

we conclude that

$$
R \leq R_{1}+R_{2} \leq A_{2}+A_{1}^{2}
$$

In the following, let $C_{1}, C_{2}, \ldots$ suitable constants. As $A_{1}^{2} \leq A_{2}$ by Jensen's inequality, and as

$$
A_{2} \leq\|\lambda\|^{2} \cdot C_{1} \cdot \int_{K} \sum_{r=1}^{n}\left|\ln \left(H_{r}(k, a)\right)\right|^{2} d k=:\|\lambda\|^{2} \cdot B_{2}
$$

we obtain $R \leq B_{2} \cdot 2\|\lambda\|^{2}$. In order to complete the proof, we must check that $B_{2}$, i.e., that the integrals

$$
\begin{equation*}
L_{r}:=\int_{K}\left|\ln \left(H_{r}(k, a)\right)\right|^{2} d k \tag{2.5}
\end{equation*}
$$

remain bounded independent of $a_{1}, \ldots, a_{n}>0$ for $r=1, \ldots, n$.
For this fix $r$. Lemma 2.1 in particular implies that for all $a_{1}, \ldots, a_{n}>0$,

$$
\Delta_{r}\left(k^{*} a k\right) \leq \sum_{1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n} a_{i_{1}} \cdot a_{i_{2}} \cdots a_{i_{r}}=\binom{n}{r} C_{r}\left(a_{1}, \ldots, a_{n}\right)
$$

and $\Delta_{r}\left(k^{*} a k\right)>0$. In other words,

$$
\begin{equation*}
0<\frac{\Delta_{r}\left(k^{*} a k\right)}{C_{r}\left(a_{1}, \ldots, a_{n}\right)}=H_{r}(k, a) \leq\binom{ n}{r} \tag{2.6}
\end{equation*}
$$

We conclude from $(2.5),(2.6)$ and Lemma 2.4 that for any $\varepsilon \in] 0,1\left[\right.$ and suitable $C_{2}=C_{2}(\varepsilon)$,

$$
L_{r} \leq C_{2} \int_{K}\left(1+H_{r}\left(a_{1}, \ldots, a_{n}\right)^{-\varepsilon}\right) d k
$$

Therefore, by Lemma 2.2,

$$
\begin{align*}
L_{r} & \leq C_{2}+C_{3} \int_{K}\left(\sum_{1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n} c_{i_{1}, \ldots, i_{r}}(k)^{-1}\right)^{\varepsilon} d k \\
& \leq C_{2}+C_{3} \cdot\binom{n}{r}^{\varepsilon} \sum_{1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n} \int_{K} c_{i_{1}, \ldots, i_{r}}(k)^{-\varepsilon} d k \tag{2.7}
\end{align*}
$$

The right hand side of (2.7) is independent of $a_{1}, \ldots, a_{n}$, and, by the definition of the $c_{i_{1}, \ldots, i_{r}}(k)$ in Lemma 2.1, $\int_{K} c_{i_{1}, \ldots, i_{r}}(k)^{-\varepsilon} d k$ is independent of $1 \leq i_{1}<i_{2}<\ldots<i_{r} \leq n$. Therefore, it suffices to check that

$$
I_{r}:=\int_{K} c_{1, \ldots, r}(k)^{-\varepsilon} d k=\int_{K} \Delta_{r}\left(k^{*}\left(\begin{array}{cc}
I_{r} & 0 \\
0 & 0
\end{array}\right) k\right)^{-\varepsilon} d k<\infty
$$

For this, we write $k$ as block matrix $k=\left(\begin{array}{cc}k_{r} & * \\ * & k_{n-r}\end{array}\right)$ with $k_{r} \in M_{r}(\mathbb{C})$ and $k_{n-r} \in M_{n-r}(\mathbb{C})$ and observe that

$$
\Delta_{r}\left(k^{*}\left(\begin{array}{cc}
I_{r} & 0 \\
0 & 0
\end{array}\right) k\right)=\Delta_{r}\left(\begin{array}{cc}
k_{r}^{*} k_{r} & * \\
* & *
\end{array}\right)=\left|\operatorname{det} k_{r}\right|^{2} .
$$

We thus have to check that $\int_{K}\left|\operatorname{det} k_{r}\right|^{-2 \varepsilon} d k<\infty$. As this is a consequence of the following lemma, the proof of the theorem is complete.
2.5 Lemma. Keep the block matrix notation above. For $\varepsilon<1 / 2$,

$$
\int_{K}\left|\operatorname{det} k_{r}\right|^{-2 \varepsilon} d k<\infty
$$

Proof. The statement is obvious for $r=n$. Moreover, by Lemma 2.3 we may assume that $1 \leq r \leq n / 2$ which we shall assume now. In this case, we introduce the matrix ball

$$
B_{r}:=\left\{w \in M_{r}(\mathbb{F}): w^{*} w \leq I_{r}\right\}
$$

as well as the ball $B:=\left\{y \in M_{1, r}(\mathbb{F}) \equiv \mathbb{F}^{n}:\|y\|_{2}^{2} \leq 1\right\}$. We conclude from the truncation lemma 2.1 of [R2] that

$$
\frac{1}{\kappa_{r}} \int_{K}\left|\operatorname{det} k_{r}\right|^{-2 \varepsilon} d k=\int_{B_{r}}|\operatorname{det} w|^{-2 \varepsilon} \Delta\left(I_{r}-w^{*} w\right)^{(n-2 r+1) \cdot d / 2-1} d w
$$

where $d w$ is the usual Lebesgue measure on the ball $B_{r}$ and

$$
\kappa_{r}:=\left(\int_{B_{r}} \operatorname{det}\left(I_{r}-w^{*} w\right)^{(n-2 r+1) \cdot d / 2-1} d w\right)^{-1}
$$

Moreover, by Lemma 3.7 and Corollary 3.8 of [R1], the mapping $P: B^{r} \rightarrow B_{r}$ with

$$
P\left(y_{1}, \ldots, y_{r}\right):=\left(\begin{array}{c}
y_{1}  \tag{2.8}\\
y_{2}\left(I_{r}-y_{1}^{*} y_{1}\right)^{1 / 2} \\
\vdots \\
y_{r}\left(I_{r}-y_{r-1}^{*} y_{r-1}\right)^{1 / 2} \cdots\left(I_{r}-y_{1}^{*} y_{1}\right)^{1 / 2}
\end{array}\right)
$$

establishes a diffeomorphism such that the image of the measure $\operatorname{det}\left(I_{r}-w^{*} w\right)^{(n-2 r+1) \cdot d / 2-1} d w$ under $P^{-1}$ is $\prod_{j=1}^{r}\left(1-\left\|y_{j}\right\|_{2}^{2}\right)^{(n-r-j+1) \cdot d / 2-1} d y_{1} \ldots d y_{r}$. Moreover, we show in Lemma 2.6 below that

$$
\operatorname{det} P\left(y_{1}, \ldots, y_{r}\right)=\operatorname{det}\left(\begin{array}{c}
y_{1} \\
\vdots \\
y_{r}
\end{array}\right)
$$

We thus conclude that

$$
\int_{K}\left|\operatorname{det} k_{r}\right|^{-2 \varepsilon} d k=\frac{1}{\kappa_{r}} \int_{B} \ldots \int_{B}\left|\operatorname{det}\left(\begin{array}{c}
y_{1}  \tag{2.9}\\
\vdots \\
y_{r}
\end{array}\right)\right|^{-2 \varepsilon} \prod_{j=1}^{r}\left(1-\left\|y_{j}\right\|_{2}^{2}\right)^{(n-r-j+1) \cdot d / 2-1} d y_{1} \ldots d y_{r}
$$

This integral is finite for $\varepsilon<1 / 2$, as one can use Fubini with an one-dimensional inner integral w.r.t. the ( 1,1 )-variable. After this inner integration, no further singularities appear from the determinant-part in the remaining integral.
2.6 Lemma. Keep the notations of the preceding proof. For all $y_{1}, \ldots, y_{n} \in B$,

$$
\operatorname{det} P\left(y_{1}, \ldots, y_{r}\right)=\operatorname{det}\left(\begin{array}{c}
y_{1} \\
\vdots \\
y_{r}
\end{array}\right) \text {. }
$$

Proof. Fix $y_{1} \in B$. The mapping $y \mapsto y\left(I_{r}-y_{1}^{*} y_{1}\right)^{1 / 2}$ on $B$ has the following form: If $y$ is written as $y=a y_{1}+y^{\perp}$ in a unique way with $a \in \mathbb{F}$ and $y^{\perp} \perp y_{1}$, then $y\left(I_{r}-y_{1}^{*} y_{1}\right)^{1 / 2}=$ $\sqrt{1-\left\|y_{1}\right\|_{2}^{2}} \cdot a y_{1}+y^{\perp}$ (write $I_{r}-y_{1}^{*} y_{1}$ in an orthonormal basis with $y_{1} /\left\|y_{1}\right\|_{2}$ as a member!). Using linearity of the determinant in all lines, we thus conclude that
$\operatorname{det}\left(\begin{array}{c}y_{1} \\ y_{2}\left(I_{r}-y_{1}^{*} y_{1}\right)^{1 / 2} \\ \vdots \\ y_{r}\left(I_{r}-y_{r-1}^{*} y_{r-1}\right)^{1 / 2} \cdots\left(I_{r}-y_{1}^{*} y_{1}\right)^{1 / 2}\end{array}\right)=\operatorname{det}\left(\begin{array}{c}y_{1} \\ y_{2} \\ y_{3}\left(I_{r}-y_{2}^{*} y_{2}\right)^{1 / 2} \\ \vdots \\ y_{r}\left(I_{r}-y_{r-1}^{*} y_{r-1}\right)^{1 / 2} \cdots\left(I_{r}-y_{2}^{*} y_{2}\right)^{1 / 2}\end{array}\right)$.
The lemma now follows by an obvious induction.

## 3 Moments and the proof of a central limit theorem

In this section we prove Theorem 1.2 and related results. We start with some facts about the moment functions of Section 1. The first result concerns an estimate for $m_{1}$.
3.1 Lemma. For $r=1, \ldots, n$ let

$$
s_{r}(g):=m_{(1,0, \ldots, 0)}(g)+\cdots+m_{(0, \ldots, 0,1,0, \ldots, 0)}(g) \quad(g \in G)
$$

be the sum of the firstr moment functions of first order. Moreover, let $\sigma_{1}(a) \geq \ldots \geq \sigma_{n}(a)>$ 0 be the ordered eigenvalues of a positive definite $n \times n$ matrix $a$. Then:
(1) $s_{n}(g)=\ln \operatorname{det}\left(g g^{*}\right)$.
(2) There is a constant $C=C(n)$ such that for all $r=1, \ldots, n$ and $g \in G$,

$$
0 \leq \ln \sigma_{1}\left(g g^{*}\right)+\ldots+\ln \sigma_{r}\left(g g^{*}\right)-s_{r}(g) \leq C .
$$

(3) There is a constant $C=C(n)$ such that for all $g \in G$

$$
\left\|2 \ln \sigma_{\text {sing }}(g)-m_{\mathbf{1}}(g)\right\| \leq C .
$$

Proof. We may assume that $g g^{*}=\operatorname{diag}\left(a_{1}, \ldots, a_{n}\right)$ with $a_{l}=\sigma_{l}\left(g g^{*}\right)(l=1 \ldots, n)$. The integral representation (1.3) implies that

$$
s_{r}(g)=\int_{K} \ln \Delta_{r}\left(k^{*} g g^{*} k\right) d k .
$$

This proves (1) and, in combination with Lemma 2.1, the first inequality in (2). For the second inequality of (2), we use the notations of Lemmas 2.1 and 2.2 . By the proof of Lemma 2.2, we have for $k \in K$,

$$
a_{1} \cdot a_{2} \ldots a_{r} \leq\binom{ n}{r} C_{r}\left(a_{1}, \ldots, a_{n}\right) \leq \max _{1 \leq i_{1}<\ldots<i_{r} \leq n} \frac{\ln \Delta_{r}\left(k^{*} g g^{*} k\right)}{c_{i_{1}, \ldots, i_{r}}(k)} .
$$

Therefore,

$$
\ln \sigma_{1}\left(g g^{*}\right)+\ldots+\ln \sigma_{r}\left(g g^{*}\right)=\int_{K} \ln \left(a_{1} \cdot a_{2} \ldots a_{r}\right) d k \leq \int_{K} \ln \Delta_{r}\left(k^{*} g g^{*} k\right) d k+M
$$

for

$$
M:=\int_{K} \max _{1 \leq i_{1}<\ldots<i_{r} \leq n} \frac{1}{c_{i_{1}, \ldots, i_{r}}(k)} d k \leq \sum_{1 \leq i_{1}<\ldots<i_{r} \leq n} \int_{K} \ln \left(c_{i_{1}, \ldots, i_{r}}(k)^{-1}\right) d k
$$

As by the definition of $c_{i_{1}, \ldots, i_{r}}(k)$ all integrals in the sum are obviously equal, it suffices to show that

$$
\int_{K} \ln \left(c_{1, \ldots, r}(k)^{-1}\right) d k=-\int_{K} \ln \Delta_{r}\left(k^{*}\left(\begin{array}{cc}
I_{r} & 0 \\
0 & 0
\end{array}\right) k\right) d k
$$

is finite. But this follows immediately from Lemma 2.5. This proves (2). Finally, (3) is a consequence of (2).

Lemma 3.1(3) implies that there exists $C=C(n)>0$ such that for all $g \in G$,

$$
\begin{equation*}
\left|e^{2 i \lambda \cdot \ln \sigma_{\text {sing }}(g)}-e^{i \lambda \cdot m_{1}(g)}\right| \leq C \cdot\|\lambda\| . \tag{3.1}
\end{equation*}
$$

Therefore, we conclude from Theorem 1.1:
3.2 Corollary. There exists a constant $C=C(n)>0$ such that for all $g \in G$,

$$
\left\|\varphi_{i \rho-\lambda}(g)-e^{2 i \lambda \cdot \ln \sigma_{s i n g}(g)}\right\| \leq C \cdot\|\lambda\| .
$$

3.3 Remark. It can be easily checked (e.g. for $n=2$ from explicit formulas in $[\mathrm{K}]$ ) that the uniform orders $\|\lambda\|^{2}$ and $\|\lambda\|$ in Theorem 1.1 and Corollary 3.2 respectively are sharp. We note that Corollary 3.2 is closely related to the Harish-Chandra expansion of the spherical functions; see e.g. Opdam [O] and Lemma I.4.2.2 of [HS] in the context of Heckman-Opdam hypergeometric functions which includes our setting. We also remark that in the proof of the CLT 1.2 below Corollary 3.2 would be sufficient instead of the stronger Theorem 1.1. On the other hand, Theorem 1.1 leads generally to stronger rates of convergence in the CLT; see e.g. Theorem 4.2 of [V2] for the rank one case.

We shall also need the following estimate which follows immediately from the integral representation (1.2):
3.4 Lemma. For all $g \in G$ and $l \in \mathbb{N}_{0}^{n}$,

$$
\left|\frac{\partial^{|l|}}{\partial \lambda^{l}} \varphi_{i \rho-\lambda}(g)\right| \leq m_{l}(g) .
$$

Let $m \in \mathbb{N}_{0}$ and $\nu \in M^{1}(G)$ a $K$-biinvariant probability measure. We say that $\nu$ admits finite $m$-th modified moments if in the notation of the introduction on the moment functions,

$$
m_{(m, 0, \ldots, 0)}, m_{(0, m, 0, \ldots, 0)}, \ldots, m_{(0, \ldots, 0, m)} \in L^{1}(G, \nu)
$$

It follows immediately from (1.3) and Hölder's inequality that in this case all moment functions of order at most $m$ are $\nu$-integrable. Moreover, this moment condition implies a corresponding differentiability of the spherical Fourier transform of $\nu$ :
3.5 Lemma. Let $m \in \mathbb{N}_{0}$ and $\nu \in M^{1}(G)$ a $K$-biinvariant probability measure with finite $m$-th moments. Then the spherical Fourier transform

$$
\tilde{\nu}: \mathbb{R}^{n} \rightarrow \mathbb{C}, \quad \lambda \mapsto \int_{G} \varphi_{i \rho-\lambda}(g) d \nu(g)
$$

is $m$-times continuously partially differentiable, and for all $l \in \mathbb{N}_{0}^{n}$ with $|l| \leq m$,

$$
\begin{equation*}
\frac{\partial^{|l|}}{\partial \lambda^{l}} \tilde{\nu}(\lambda)=\int_{G} \frac{\partial^{|l|}}{\partial \lambda^{l}} \varphi_{i \rho-\lambda}(g) d \nu(g) \tag{3.2}
\end{equation*}
$$

In particular,

$$
\begin{equation*}
\frac{\partial^{|l|}}{\partial \lambda^{l}} \tilde{\nu}(0)=(-i)^{|l|} \int_{G} m_{l}(g) d \nu(g) \tag{3.3}
\end{equation*}
$$

Proof. We proceed by induction: The case $m=0$ is trivial, and for $m \rightarrow m+1$ we observe that by our assumption all moments of lower order exist, i.e., (3.2) is available for all $|l| \leq m$. It now follows from Lemma 3.4 and the well-known result about parameter integrals that a further partial derivative and the integration can be interchanged. Finally, (3.3) follows from (3.2) and (1.3). Continuity of the derivatives is also clear by Lemma 3.4.

We next turn to the positive (semi)definiteness of the modified covariance matrix $\sigma^{2}(\nu)$ for biinvariant measures with finite second modified moments. We start with measures concentrated on a double coset:
3.6 Lemma. Let $n \geq 2, g \in G$, and $\Sigma^{2}(g):=m_{\mathbf{2}}(g)-m_{\mathbf{1}}(g)^{t} m_{\mathbf{1}}(g)$.
(1) $\Sigma^{2}(g)$ is positive semidefinite.
(2) If $g g^{*}$ is not a multiple of the identity matrix, then $\Sigma^{2}(g)$ has rank $n-1$.
(3) If $g g^{*}$ is a multiple of the identity matrix, then $\Sigma^{2}(g)=0$.

Proof. Let $a_{1}, \ldots, a_{n} \in \mathbb{R}$ with $a_{1}^{2}+\ldots+a_{n}^{2}>0$ and the row vector $a=\left(a_{1}, \ldots, a_{n}\right)$. Put

$$
f_{1}(k):=\ln \Delta_{1}\left(k^{*} g g^{*} k\right) \quad \text { and } \quad f_{l}(k):=\ln \Delta_{l}\left(k^{*} g g^{*} k\right)-\ln \Delta_{l-1}\left(k^{*} g g^{*} k\right) \quad(l=2, \cdots, n)
$$

Then, by (1.3), (1.4), (1.5), and the Cauchy-Schwarz inequality,

$$
a\left(m_{\mathbf{2}}(g)-m_{\mathbf{1}}(g)^{t} m_{\mathbf{1}}(g)\right) a^{t}=\int_{K}\left(\sum_{l=1}^{n} a_{l} f_{l}(k)\right)^{2} d k-\left(\int_{K} \sum_{l=1}^{n} a_{l} f_{l}(k) d k\right)^{2} \geq 0
$$

Moreover, this expression is equal to 0 if and only if the function
$k \mapsto \sum_{l=1}^{n} a_{l} f_{l}(k)=\left(a_{1}-a_{2}\right) \ln \Delta_{1}\left(k^{*} g g^{*} k\right)+\cdots+\left(a_{n-1}-a_{n}\right) \ln \Delta_{n-1}\left(k^{*} g g^{*} k\right)+a_{n} \ln \Delta_{n}\left(k^{*} g g^{*} k\right)$
is constant on $K$. As $k \mapsto \ln \Delta_{n}\left(k^{*} g g^{*} k\right)$ is constant on $K$, and as under the condition of (2), the functions $k \mapsto \ln \Delta_{r}\left(k^{*} g g^{*} k\right)(r=1, \ldots, n-1)$ and the constant function 1 are linearly independent on $K$ by Corollary 4.2 in the appendix, the function $k \mapsto \sum_{l=1}^{n} a_{l} f_{l}(k)$ is constant on $K$ precisely for $a_{1}=a_{2}=\ldots=a_{n}$. This proves (2). Part (3) is obvious.

The arguments of the preceding proof lead to the following characterization of $K$-biinvariant measures with positive definite covariance matrices:
3.7 Lemma. Let $\nu \in M^{1}(G)$ be a $K$-biinvariant probability measure having second modified moments. Then $\Sigma^{2}(\nu)$ is positive definite if and only if supp $\nu$ is not contained in the subgroups $\left\{c I_{n}: c \in \mathbb{F}, c \neq 0\right\}$ and $S L_{n}(\mathbb{F})$.

We now turn to the proof of the CLT:
Proof of Theorem 1.2. Let $\nu \in M^{1}(G)$ be a $K$-biinvariant probability measure with finite second modified moments. Let $\left(X_{k}\right)_{k \geq 1}$ be i.i.d. $G$-valued random variables with distribution $\nu$ and $S_{k}:=X_{1} \cdot X_{2} \cdots X_{k}$. Let $\lambda \in \mathbb{R}^{n}$. As the functions $\varphi_{i \rho-\lambda}$ are bounded on $G$ (by the integral representation (1.2)) and multiplicative w.r.t. $K$-biinvariant measures, we have

$$
E\left(\varphi_{i \rho-\lambda / \sqrt{k}}\left(S_{k}\right)\right)=\int_{G} \varphi_{i \rho-\lambda / \sqrt{k}}(g) d \nu^{(k)}(g)=\left(\int_{G} \varphi_{i \rho-\lambda / \sqrt{k}}(g) d \nu(g)\right)^{k}=\tilde{\nu}(\lambda / \sqrt{k})^{k} .
$$

We now use Taylor's formula, Lemma 3.5, and

$$
m_{\mathbf{2}}(\nu):=\int_{G} m_{\mathbf{2}}(g) d \nu(g)=\Sigma^{2}(\nu)+m_{\mathbf{1}}(\nu)^{t} m_{\mathbf{1}}(\nu)
$$

and obtain

$$
\begin{aligned}
& \exp \left(i \lambda \cdot m_{\mathbf{1}}(\nu) \sqrt{k}\right) \cdot E\left(\varphi_{i \rho-\lambda / \sqrt{k}}\left(S_{k}\right)\right)=\left(\exp \left(i \lambda \cdot m_{\mathbf{1}}(\nu) / \sqrt{k}\right) \cdot \tilde{\nu}(\lambda / \sqrt{k})\right)^{k} \\
&=\left(\left[1+\frac{i \lambda \cdot m_{\mathbf{1}}(\nu)}{\sqrt{k}}-\frac{\left(\lambda \cdot m_{\mathbf{1}}(\nu)\right)^{2}}{2 k}+o\left(\frac{1}{k}\right)\right] \cdot\left[1-\frac{i \lambda \cdot m_{\mathbf{1}}(\nu)}{\sqrt{k}}-\frac{\lambda m_{\mathbf{2}}(\nu) \lambda^{t}}{2 k}+o\left(\frac{1}{k}\right)\right]\right)^{k} \\
&=\left(\left[1+\frac{i \lambda \cdot m_{\mathbf{1}}(\nu)}{\sqrt{k}}-\frac{\left(\lambda \cdot m_{\mathbf{1}}(\nu)\right)^{2}}{2 k}+o\left(\frac{1}{k}\right)\right] .\right. \\
&\left.\quad \times\left[1-\frac{i \lambda \cdot m_{\mathbf{1}}(\nu)}{\sqrt{k}}-\frac{\lambda\left(\Sigma^{2}(\nu)+m_{\mathbf{1}}(\nu)^{t} m_{\mathbf{1}}(\nu)\right) \lambda^{t}}{2 k}+o\left(\frac{1}{k}\right)\right]\right)^{k} \\
&=\left(1-\frac{\lambda \Sigma^{2}(\nu) \lambda^{t}}{2 k}+o\left(\frac{1}{k}\right)\right)^{k} .
\end{aligned}
$$

Therefore,

$$
\lim _{k \rightarrow \infty} \exp \left(i \lambda \cdot m_{1}(\nu) \sqrt{k}\right) \cdot E\left(\varphi_{i \rho-\lambda / \sqrt{k}}\left(S_{k}\right)\right)=\exp \left(-\lambda \Sigma^{2}(\nu) \lambda^{t} / 2\right)
$$

Moreover, by Theorem 1.1,

$$
\lim _{k \rightarrow \infty} E\left(\varphi_{i \rho-\lambda / \sqrt{k}}\left(S_{k}\right)-\exp \left(-i \lambda \cdot m_{\mathbf{1}}\left(S_{k}\right) / \sqrt{k}\right)\right)=0
$$

We conclude that

$$
\lim _{k \rightarrow \infty} \exp \left(-i \lambda \cdot\left(m_{\mathbf{1}}\left(S_{k}\right)-k \cdot m_{\mathbf{1}}(\nu)\right) / \sqrt{k}\right)=\exp \left(-\lambda \Sigma^{2}(\nu) \lambda^{t} / 2\right)
$$

for all $\lambda \in \mathbb{R}^{n}$. Levy's continuity theorem for the classical $n$-dimensional Fourier transform now implies that $\left(m_{1}\left(S_{k}\right)-k \cdot m_{1}(\nu)\right) / \sqrt{k}$ tends in distribution to $N\left(0, \Sigma^{2}(\nu)\right)$. By the estimate of Lemma 3.1(1), this immediately implies Theorem 1.2.

On the basis of Theorem 1.1, also a Berry-Esseen-type estimate with the order $O\left(k^{-1 / 3}\right)$ of convergence can be derived. As the details are technical, but quite similar to the proof of the corresponding rank-one-case in Theorem 4.2 of [V2], we here omit details. We also mention that Theorem 1.1 can be also used to derive further CLTs e.g. with stable distributions with domains of attraction or a Lindeberg-Feller CLT. The details of proof then would be also very similar to the classical cases for sums of iid random variables.

## 4 Appendix

Here we collect some results from linear algebra which are needed in Section 3.
4.1 Lemma. Let $x_{1}, \ldots, x_{n} \in \mathbb{R}$. Then

$$
\begin{aligned}
& \operatorname{det}\left(\begin{array}{cccccc}
x_{1} & x_{2} & x_{3} & x_{4} & \cdots & x_{n} \\
x_{1}+x_{2} & x_{2}+x_{1} & x_{3}+x_{2} & x_{4}+x_{2} & \cdots & x_{n}+x_{2} \\
x_{1}+x_{2}+x_{3} & x_{2}+x_{1}+x_{3} & x_{3}+x_{2}+x_{1} & x_{4}+x_{2}+x_{3} & \cdots & x_{n}+x_{2}+x_{3} \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\
\sum_{l=1}^{n} x_{l} & \sum_{l=1}^{n} x_{l} & \sum_{l=1}^{n} x_{l} & \sum_{l=1}^{n} x_{l} & \cdots & \sum_{l=1}^{n} x_{l}
\end{array}\right)= \\
& =\left(x_{1}+x_{2}+\cdots+x_{n}\right) \cdot\left(x_{1}-x_{2}\right) \cdot\left(x_{1}-x_{3}\right) \cdots\left(x_{1}-x_{n}\right)
\end{aligned}
$$

Proof. The determinant is a homogeneous polynomial in the the variables $x_{1}, \ldots, x_{n}$ of degree $n$. Moreover, the monomial $x_{1}^{n}$ appears in this polynomial with coefficient 1 , and for given $x_{2}, \ldots, x_{n}$, the determinant is a polynomial in the variable $x_{1}$ where $-\left(x_{2}+\cdots+x_{n}\right), x_{2}$, $x_{3}, \ldots, x_{n}$ are the zeros of this polynomial. This leads readily to the claim.
4.2 Corollary. Let $a_{1}, \ldots, a_{n}>0$ numbers such that at least two of them are different. Consider the diagonal matrix $a=\operatorname{diag}\left(a_{1}, \ldots, a_{n}\right)$. Then the functions $k \mapsto \ln \Delta_{r}\left(k^{*} a k\right)$ with $r=1, \ldots, n-1$ and the constant function 1 on $K=U_{n}(\mathbb{F})$ are linearly independent.

Proof. Without loss of generality, $a_{1}$ is different from $a_{2}, \ldots, a_{n}$. Now consider the $n$ permutation matrices $k_{l}$ which permute the rows 1 and $l$ and leave the other rows invariant for $l=$ $1, \cdots, n$. Then, using the notation $x_{l}:=\ln a_{l}$, the number $\ln \Delta_{r}\left(k_{j}^{*} a k_{j}\right)$ is precisely the $r, l-$ entry of the matrix in Lemma 4.1. Therefore, by Lemma 4.1, $\operatorname{det}\left(\left(\ln \Delta_{r}\left(k_{j}^{*} a k_{j}\right)\right)_{r, j=1, \ldots, n}\right) \neq 0$ for $x_{1}+\ldots+x_{n} \neq 0$, i.e., $a$ with $\operatorname{det} a \neq 1$. As $\ln \Delta_{n}\left(k^{*} a k\right)$ is constant, this proves the statement of the corollary for $\operatorname{det} a \neq 1$. The case $\operatorname{det} a=1$ can be easily derived by considering $2 a$ instead of $a$ in the preceding argument.

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