# Riesz measures and Wishart laws associated to quadratic maps 

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#### Abstract

We introduce a natural definition of Riesz measures and Wishart laws associated to an $\Omega$-positive (virtual) quadratic map, where $\Omega \subset \mathbb{R}^{n}$ is a regular open convex cone. We give a general formula for moments of the Wishart laws. Moreover, if the quadratic map has an equivariance property under the action of a linear group acting on the cone $\Omega$ transitively, then the associated Riesz measure and Wishart law are described explicitly by making use of theory of relatively invariant distributions on homogeneous cones.


## §1. Introduction.

Riesz measures and distributions on convex cones form one of fundamental tools of harmonic analysis and of the theory of the wave equation, cf. [4] in the case of symmetric cones and [6, [11] for homogeneous cones. Moreover, exponential families generated by Riesz measures are composed of Wishart laws and are of great significance in random matrix theory and in statistics.

Wishart laws are probability distributions on symmetric or Hermitian matrices with very important applications in multivariate statistics. Their role in statistics is due to two reasons:

- they are probability distributions of the maximum likelihood estimator(MLE) of the covariance matrix in a multivariate normal sample ([20, [1]).
- in Bayesian statistics, Wishart laws form a Diaconis-Ylvisaker family ([3]) of prior distributions for the covariance parameter in a covariance selection model (19]).

On the other hand, recent developments in random matrix theory of chiral Gaussian ensembles containing Wishart laws, are intense and motivated by applications

[^0]in mathematical physics, cf. [16] and references therein.

These numerous modern applications of Wishart laws make it necessary to develop the theory of Wishart laws and Riesz measures on more general cones than in the classical case of the symmetric cones of real symmetric or complex Hermitian matrices. For example, in an $r$-dimensional Gaussian model $X$, if the marginal variables $X_{i}$ and $X_{j}$ are known to be conditionally independent given all the other variables, the statistical analysis of the covariance matrix of $X$ must be done on the cone $\mathcal{P}$ of positive definite symmetric matrices $Y$ with $Y_{i j}=Y_{j i}=0$ and on its dual cone $\mathcal{Q}$ ([19]). The cones $\mathcal{P}$ and $\mathcal{Q}$ are usually no longer symmetric. This led to some important papers in recent statistical and probabilistic literature about Wishart laws on more general cones: homogeneous cones ([1) or cones related to graphical models (19). In these papers, Wishart laws are introduced via their density functions (see Section 3.8).

In our paper we construct and study Riesz measures and Wishart laws on regular convex cones via quadratic maps. For a regular open convex cone $\Omega \subset \mathbb{R}^{n}$ and an $\Omega$-positive quadratic map $q: \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$, the Riesz measure associated to $q$ is defined as the image of the Lebesgue measure $d x$ on $\mathbb{R}^{m}$ by $q$. Wishart laws studied in this paper are obtained from $\mathbb{R}^{m}$-valued normal random vectors $X$ as the law of $Y:=q(X) / 2$. This is a classical natural approach to Riesz measures (4) and Wishart laws $([20,4])$ and we propose to extend it to a much more general setting.

In Section 2 of the paper we explain the details of the quadratic construction of Riesz measures on regular convex cones and next we define the corresponding Wishart laws. We compute their Laplace transforms, what is the starting point to get formulas for their expectation, covariance and higher moments (Theorems 2.5 and (2.6). More general Riesz and Wishart distributions associated to virtual quadratic maps are introduced in Section 2.3. Moments formulas are generalized (Theorem [2.9). Group equivariance of Wishart laws is studied at the end of the section.

Section 3 of the article is thoroughly devoted to the case when $\Omega$ is a homogeneous cone and the quadratic map $q$ is homogeneous. A crucial role in the analysis of these maps and of related Riesz measures and Wishart laws is played by a matrix realization of any homogeneous cone, coming from [14] and explained in Section 3.2.

It allows, among others, to define basic and standard quadratic maps in Sections 3.3 and 3.4. They play a role of generators for homogeneous quadratic maps $q$ needed to construct all Riesz measures and Wishart laws on $\Omega$. Next we apply the results of [11] on Gindikin-Riesz distributions on $\Omega$ and on the orbit decomposition of $\bar{\Omega}$, the closure of $\Omega$. We explain the relation between Riesz measures related to homogeneous quadratic maps and the Gindikin-Riesz distributions on $\Omega$ (Theorem 3.9). In Section 3.7, we prove the Bartlett decomposition for the Wishart laws on homogeneous cones (Theorems 3.11 and 3.13).

Families of Wishart laws that we construct and study in Section 3 comprise Wishart distributions studied in papers [1] and [19] (homogeneous case) and are significantly bigger: we describe all singular Wishart laws and many more absolutely continuous Wishart laws than in papers [1] and [19]. For the symmetric cone case, our Wishart laws cover the ones studied in 10 as well. All the results of Section 2 apply to them, in particular the formulas for the moments.

Throughout the paper and from its very beginning all our concepts are illustrated on important examples, including a non-homogeneous cone (Example 1 Section 2), symmetric cones of positive definite real symmetric matrices $\Pi_{r}$ (Example 2 Section 2), 4-dimensional Lorentz cone(Example 5 Section 3), and non-symmetric but homogeneous Vinberg cone and its dual(Example 3 Section 2).

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## §2. Riesz measure and Wishart law on a convex cone.

2.1. Regular cones and quadratic maps. In this paper, an open convex cone $\Omega \subset \mathbb{R}^{n}$ is always assumed to be regular, that is, $\bar{\Omega} \cap(-\bar{\Omega})=\{0\}$, where $\bar{\Omega}$ denotes the closure of $\Omega$. Then the dual cone $\Omega^{*}:=\left\{\eta \in\left(\mathbb{R}^{n}\right)^{*} ;\langle y, \eta\rangle>0(\forall y \in \bar{\Omega} \backslash\{0\})\right\}$ is a regular open convex cone again in the dual vector space $\left(\mathbb{R}^{n}\right)^{*}$, and we have $\left(\Omega^{*}\right)^{*}=\Omega$. An $\mathbb{R}^{n}$-valued quadratic map $q: \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$ is said to be $\Omega$-positive if (i) $q(x) \in \bar{\Omega}$ for all $x \in \mathbb{R}^{m}$, and (ii) $q(x)=0$ implies $x=0$. These (i) and (ii) are restated in a single condition $q(x) \in \bar{\Omega} \backslash\{0\}\left(\forall x \in \mathbb{R}^{m} \backslash\{0\}\right)$. For the quadratic map $q$, we define the associated linear $\operatorname{map} \phi=\phi_{q}:\left(\mathbb{R}^{n}\right)^{*} \rightarrow \operatorname{Sym}(m, \mathbb{R})$ in such a way that

$$
{ }^{\mathrm{t}} x \phi(\eta) x=\langle q(x), \eta\rangle \quad\left(\eta \in\left(\mathbb{R}^{n}\right)^{*}, x \in \mathbb{R}^{m}\right) .
$$

Then the $\Omega$-positivity of $q$ is equivalent to the following property of $\phi$ :

$$
\begin{equation*}
\eta \in \Omega^{*} \Rightarrow \phi(\eta) \text { is positive definite. } \tag{2.1}
\end{equation*}
$$

Example 1. Let $\Omega$ be the open convex cone in $\mathbb{R}^{3}$ defined by

$$
\begin{align*}
\Omega & :=\left\{t_{1}\left(\begin{array}{l}
0 \\
0 \\
1
\end{array}\right)+t_{2}\left(\begin{array}{l}
1 \\
0 \\
1
\end{array}\right)+t_{3}\left(\begin{array}{l}
1 \\
1 \\
1
\end{array}\right)+t_{4}\left(\begin{array}{l}
0 \\
1 \\
1
\end{array}\right) ; t_{1}, t_{2}, t_{3}, t_{4}>0\right\} \\
& =\left\{\left(\begin{array}{l}
y_{1} \\
y_{2} \\
y_{3}
\end{array}\right) \in \mathbb{R}^{3} ; y_{1}>0, y_{2}>0,-y_{1}+y_{3}>0,-y_{2}+y_{3}>0\right\} . \tag{2.2}
\end{align*}
$$

If we identify $\left(\mathbb{R}^{3}\right)^{*}$ with $\mathbb{R}^{3}$ by $\langle y, \eta\rangle:=y_{1} \eta_{1}+y_{2} \eta_{2}+y_{3} \eta_{3}\left(y, \eta \in \mathbb{R}^{3}\right)$, we have

$$
\begin{aligned}
\Omega^{*} & =\left\{t_{1}\left(\begin{array}{l}
1 \\
0 \\
0
\end{array}\right)+t_{2}\left(\begin{array}{l}
0 \\
1 \\
0
\end{array}\right)+t_{3}\left(\begin{array}{c}
-1 \\
0 \\
1
\end{array}\right)+t_{4}\left(\begin{array}{c}
0 \\
-1 \\
1
\end{array}\right) ; t_{1}, t_{2}, t_{3}, t_{4}>0\right\} \\
& =\left\{\left(\begin{array}{l}
\eta_{1} \\
\eta_{2} \\
\eta_{3}
\end{array}\right) \in \mathbb{R}^{3} ; \eta_{3}>0, \eta_{1}+\eta_{3}>0, \eta_{1}+\eta_{2}+\eta_{3}>0, \eta_{2}+\eta_{3}>0\right\},
\end{aligned}
$$

see [13]. Let $q: \mathbb{R}^{4} \rightarrow \mathbb{R}^{3}$ be the quadratic map given by

$$
q(x):=\left(x_{1}\right)^{2}\left(\begin{array}{l}
0 \\
0 \\
1
\end{array}\right)+\left(x_{2}\right)^{2}\left(\begin{array}{l}
1 \\
0 \\
1
\end{array}\right)+\left(x_{3}\right)^{2}\left(\begin{array}{l}
1 \\
1 \\
1
\end{array}\right)+\left(x_{4}\right)^{2}\left(\begin{array}{l}
0 \\
1 \\
1
\end{array}\right) \quad\left(x \in \mathbb{R}^{4}\right)
$$

Clearly, this $q$ is $\Omega$-positive. By a simple calculation, we have

$$
\phi(\eta)=\left(\begin{array}{cccc}
\eta_{3} & 0 & 0 & 0 \\
0 & \eta_{1}+\eta_{3} & 0 & 0 \\
0 & 0 & \eta_{1}+\eta_{2}+\eta_{3} & 0 \\
0 & 0 & 0 & \eta_{2}+\eta_{3}
\end{array}\right) \quad\left(\eta \in\left(\mathbb{R}^{3}\right)^{*}\right) .
$$

Example 2. Let $\Pi_{r}$ be the set of positive definite real symmetric matrices of size $r$. Then $\Pi_{r}$ is a regular open convex cone in the vector space $\operatorname{Sym}(r, \mathbb{R})$ of real symmetric matrices. If we identify the space $\operatorname{Sym}(r, \mathbb{R})$ with its dual vector space by the inner product $\langle y, \eta\rangle:=\operatorname{tr}(y \eta)(y, \eta \in \operatorname{Sym}(r, \mathbb{R}))$, then the dual cone $\Pi_{r}^{*}$ coincides with $\Pi_{r}$. We define $q_{r, s}: \operatorname{Mat}(r, s ; \mathbb{R}) \rightarrow \operatorname{Sym}(r, \mathbb{R})$ by

$$
q_{r, s}(x)=x^{\mathrm{t}} x \quad(x \in \operatorname{Mat}(r, s ; \mathbb{R}))
$$

Then $q_{r, s}$ is $\Pi_{r}$-positive. We denote the $(i, j)$ component of $x \in \operatorname{Mat}(r, s ; \mathbb{R})$ by $x_{r(j-1)+i}$, so that $\operatorname{Mat}(r, s ; \mathbb{R})$ is identified with $\mathbb{R}^{r s}$. Then we have for $\eta \in \operatorname{Sym}(r, \mathbb{R})$

$$
\phi(\eta)=\left(\begin{array}{cccc}
\eta & & & \\
& \eta & & \\
& & \ddots & \\
& & & \eta
\end{array}\right) \in \operatorname{Sym}(r s, \mathbb{R})
$$

where $\eta$ is put $s$ times. In this case, the map $\phi: \operatorname{Sym}(r, \mathbb{R}) \rightarrow \operatorname{Sym}(r s, \mathbb{R})$ is a Jordan algebra representation, and $q$ is exactly the quadratic map associated to the representation ([4, Chapter IV, Section 4]).
Example 3. Let $\mathcal{Z}$ be a subspace of $\operatorname{Sym}(r, \mathbb{R})$, and put $\mathcal{P}:=\mathcal{Z} \cap \Pi_{r}$. Then $\mathcal{P}$ is a regular open convex cone in $\mathcal{Z}$. Let $\mathcal{Q} \subset \mathcal{Z}^{*}$ be the dual cone of $\mathcal{P}$. We shall construct a $\mathcal{Q}$-positive quadratic map $q_{\mathcal{Z}}: \mathbb{R}^{r} \rightarrow \mathcal{Z}^{*}$ whose associated linear map $\phi_{\mathcal{Z}}: \mathcal{Z} \rightarrow \operatorname{Sym}(r, \mathbb{R})$ equals the inclusion map. Let us define the surjective linear $\operatorname{map} \pi_{\mathcal{Z}^{*}}: \operatorname{Sym}(r, \mathbb{R}) \rightarrow \mathcal{Z}^{*}$ by

$$
\left\langle y, \pi_{\mathcal{Z}^{*}}(S)\right\rangle:=\operatorname{tr} y S \quad(y \in \mathcal{Z}, S \in \operatorname{Sym}(r, \mathbb{R}))
$$

Then the quadratic map $q_{\mathcal{Z}}: \mathbb{R}^{r} \rightarrow \mathcal{Z}^{*}$ is given by $q_{\mathcal{Z}}(x):=\pi_{\mathcal{Z}^{*}}\left(x^{\mathrm{t}} x\right)\left(x \in \mathbb{R}^{r}\right)$. In fact, for $x \in \mathbb{R}^{r} \backslash\{0\}$ and $y \in \mathcal{P}$ we have

$$
\begin{equation*}
\left\langle y, q_{\mathcal{Z}}(x)\right\rangle=\operatorname{tr}\left(y x^{\mathrm{t}} x\right)={ }^{\mathrm{t}} x y x>0 \tag{2.3}
\end{equation*}
$$

because $y$ is positive definite. Therefore we get $q_{\mathcal{Z}}(x) \in \overline{\mathcal{Q}} \backslash\{0\}$, so that $q_{\mathcal{Z}}$ is $\mathcal{Q}$ positive. Keeping the natural isomorphism $\left(\mathcal{Z}^{*}\right)^{*} \simeq \mathcal{Z}$ in mind, we see from (2.3) that $\phi_{\mathcal{Z}}(y)=y(y \in \mathcal{Z})$. Soon later, we shall consider the cases

$$
\mathcal{Z}:=\left\{\left(\begin{array}{ccc}
y_{11} & 0 & 0  \tag{2.4}\\
0 & y_{22} & y_{32} \\
0 & y_{32} & y_{33}
\end{array}\right) \in \operatorname{Sym}(3, \mathbb{R}) ; y_{11}, y_{22}, y_{32}, y_{33} \in \mathbb{R}\right\}
$$

and

$$
\mathcal{Z}:=\left\{\left(\begin{array}{ccc}
y_{11} & 0 & y_{31}  \tag{2.5}\\
0 & y_{22} & y_{32} \\
y_{31} & y_{32} & y_{33}
\end{array}\right) \in \operatorname{Sym}(3, \mathbb{R}) ; y_{11}, y_{21}, y_{31}, y_{22}, y_{33} \in \mathbb{R}\right\}
$$

as concrete examples. Actually, in the latter case (2.5), the cones $\mathcal{Q}$ and $\mathcal{P}$ are called the Vinberg cone and the dual Vinberg cone respectively, which are the lowest dimensional non-symmetric homogeneous cones ([23]). We shall see another realization of the Vinberg cone $\mathcal{Q}$ in (3.6) and the last paragraph of Section 3.3.

Let $I=\left\{i_{1}, i_{2}, \ldots, i_{k}\right\}$ be a subset of $\{1, \ldots, r\}$ with $1 \leq i_{1}<i_{2}<\cdots<i_{k} \leq r$, and define

$$
\begin{equation*}
R^{I}:=\left\{x \in \mathbb{R}^{r} ; x_{i}=0 \text { if } i \notin I\right\} . \tag{2.6}
\end{equation*}
$$

We denote by $q_{\mathcal{Z}}^{I}$ the restriction of $q_{\mathcal{Z}}$ to the space $R^{I} \subset \mathbb{R}^{r}$. Clearly $q_{\mathcal{Z}}^{I}: R^{I} \rightarrow \mathcal{Z}^{*}$ is $\mathcal{Q}$-positive. The associated linear map $\phi_{q_{\mathcal{Z}}^{I}}: \mathcal{Z} \rightarrow \operatorname{Sym}(k, \mathbb{R})$ gives a submatrix of elements $y \in \mathcal{Z}$, that is, $\phi_{q_{\mathcal{Z}}^{I}}(y)=\left(y_{i_{\alpha} i_{\beta}}\right)$, which we denote by $y_{I}$.

### 2.2. Riesz measures and Wishart laws associated to quadratic maps. For

 a regular open convex cone $\Omega \subset \mathbb{R}^{n}$ and an $\Omega$-positive quadratic map $q: \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$, let $\mu_{q}$ be the image of the Lebesgue measure $d x$ on $\mathbb{R}^{m}$ by $q$. Namely, the measure $\mu_{q}$ on $\mathbb{R}^{n}$ is defined in such a way that$$
\begin{equation*}
\int_{\mathbb{R}^{n}} f(y) \mu_{q}(d y)=\int_{\mathbb{R}^{m}} f(q(x)) d x \tag{2.7}
\end{equation*}
$$

for a measurable function $f$ on $\mathbb{R}^{n}$. The $\Omega$-positivity of $q$ implies that the support of $\mu_{q}$ is contained in the closure $\bar{\Omega}$ of the cone $\Omega$. By analogy to [4, Proposition VII.2.4], we call $\mu_{q}$ the Riesz measure associated to $q$.

Lemma 2.1. Let $\phi:\left(\mathbb{R}^{n}\right)^{*} \rightarrow \operatorname{Sym}(m, \mathbb{R})$ be the linear map associated to $q$. Then, for $\eta \in \Omega^{*}$, the Laplace transform $L_{\mu_{q}}(-\eta):=\int_{\mathbb{R}^{n}} e^{-\langle\eta, y\rangle} \mu_{q}(d y)$ of $\mu_{q}$ equals $\pi^{m / 2}(\operatorname{det} \phi(\eta))^{-1 / 2}$.

Proof. By definition, we have $L_{\mu_{q}}(-\eta)=\int_{\mathbb{R}^{m}} e^{-\mathrm{t} x \phi(\eta) x} d x$. Since $\phi(\eta)$ is positive definite, the assertion follows from a formula of the Gaussian integral.

Definition 2.2. The members of the exponential family $\left\{\gamma_{q, \theta}\right\}_{\theta \in-\Omega^{*}}$ generated by $\mu_{q}$ are called the Wishart laws on $\Omega$ associated to $q$. Namely,

$$
\begin{equation*}
\gamma_{q, \theta}(d y):=\frac{e^{\langle y, \theta\rangle}}{L_{\mu_{q}}(\theta)} \mu_{q}(d y) \quad\left(y \in \mathbb{R}^{n}\right) \tag{2.8}
\end{equation*}
$$

Remark 2.3. By (2.7), (2.8) and Lemma 2.1, we have for a measurable function $f$ on $\mathbb{R}^{n}$

$$
\int_{\mathbb{R}^{n}} f(y) \gamma_{q, \theta}(d y)=\pi^{-m / 2}(\operatorname{det} \phi(-\theta))^{1 / 2} \int_{\mathbb{R}^{m}} f(q(x)) e^{-\mathrm{t} x \phi(-\theta) x} d x
$$

Putting $\Sigma:=\phi(-\theta)^{-1}$ and replacing the variable $x$ by $x / \sqrt{2}$, we rewrite the righthand side as

$$
(2 \pi)^{-m / 2}(\operatorname{det} \Sigma)^{-1 / 2} \int_{\mathbb{R}^{m}} f(q(x) / 2) e^{-\mathrm{t} x \Sigma^{-1} x / 2} d x
$$

Therefore, if $X$ is an $\mathbb{R}^{m}$-valued random variable with the normal law $N\left(0, \phi(-\theta)^{-1}\right)$, then $\gamma_{q, \theta}$ is nothing else but the law of $Y:=q(X) / 2$. In particular, the classical Wishart law as defined in [20, Definition 3.1.3] coincides with our $\gamma_{q, \theta}$ in Example 2.

Proposition 2.4. Let $Y$ be an $\mathbb{R}^{n}$-valued random variable with the Wishart law $\gamma_{q, \theta}$. Then the Laplace transform $L_{\gamma_{q, \theta}}(\eta)=E\left(e^{\langle Y, \eta\rangle}\right)$ of $\gamma_{q, \theta}$ is given by

$$
L_{\gamma_{q, \theta}}(\eta)=\operatorname{det}\left(I_{m}+\phi(-\theta)^{-1} \phi(-\eta)\right)^{-1 / 2}
$$

for $\eta \in-\theta-\Omega^{*}$.
Proof. By definition, we have $L_{\gamma_{q, \theta}}(\eta)=L_{\mu_{q}}(\theta)^{-1} L_{\mu_{q}}(\eta+\theta)$. Thus the formula follows from Lemma 2.1 and the observation that $(\operatorname{det} \phi(-\theta))^{-1} \operatorname{det} \phi(-\eta-\theta)=$ $\operatorname{det}\left(\phi(-\theta)^{-1}(\phi(-\theta)+\phi(-\eta))\right)=\operatorname{det}\left(I_{m}+\phi(-\theta)^{-1} \phi(-\eta)\right)$.

We shall consider the mean and the covariance of the Wishart law $\gamma_{q, \theta}$. First we fix a notation for the directional derivative of a function, and recall some basic formulas. For a (vector-valued) function $f$ on a domain $U \subset \mathbb{R}^{N}$, a point $a \in U$ and a vector $v \in \mathbb{R}^{N}$, we denote by $D_{v} f(a)$ the directional derivative of $f$ at $a$ given by $D_{v} f(a):=\left(\frac{d}{d h}\right)_{h=0} f(a+h v)$. Now we consider some functions and their derivatives on the domain $-\Pi_{m}$ in the vector space $\operatorname{Sym}(m, \mathbb{R})$. First we set $f(x)=$ $\log \operatorname{det}(-x) \in \mathbb{R}$ for $x \in-\Pi_{m}$. Then we have

$$
\begin{equation*}
D_{v} f(a)=-\operatorname{tr}(-a)^{-1} v \quad\left(a \in-\Pi_{m}, v \in \operatorname{Sym}(m, \mathbb{R})\right) \tag{2.9}
\end{equation*}
$$

Next we observe the case $f(x)=(-x)^{-1} \in \operatorname{Sym}(m, \mathbb{R})$ for $x \in-\Pi_{m}$. Then

$$
\begin{equation*}
D_{v} f(a)=(-a)^{-1} v(-a)^{-1} \quad\left(a \in-\Pi_{m}, v \in \operatorname{Sym}(m, \mathbb{R})\right) \tag{2.10}
\end{equation*}
$$

Theorem 2.5. Let $Y$ be an $\mathbb{R}^{n}$-valued random variable with the Wishart law $\gamma_{q, \theta}$. (i) For $\eta \in\left(\mathbb{R}^{n}\right)^{*}$, one has

$$
E(\langle Y, \eta\rangle)=\operatorname{tr} \phi(-\theta)^{-1} \phi(\eta) / 2 .
$$

(ii) For $\eta, \eta^{\prime} \in\left(\mathbb{R}^{n}\right)^{*}$, one has

$$
E\left((\langle Y, \eta\rangle-M)\left(\left\langle Y, \eta^{\prime}\right\rangle-M^{\prime}\right)\right)=\operatorname{tr} \phi(-\theta)^{-1} \phi(\eta) \phi(-\theta)^{-1} \phi\left(\eta^{\prime}\right) / 2
$$

where $M:=E(\langle Y, \eta\rangle)$ and $M^{\prime}:=E\left(\left\langle Y, \eta^{\prime}\right\rangle\right)$.

Proof. Since $\left\{\gamma_{q, \theta}\right\}_{\theta \in-\Omega^{*}}$ is the exponential family generated by $\mu_{q}$, it is well known ([18]) that the mean $E(\langle Y, \eta\rangle)$ is given by the derivative $D_{\eta} \log L_{\mu}(\theta)$, while the covariance $E\left((\langle Y, \eta\rangle-M)\left(\left\langle Y, \eta^{\prime}\right\rangle-M^{\prime}\right)\right)$ equals $D_{\eta} D_{\eta^{\prime}} \log L_{\mu}(\theta)$. Thus the formulas follow from Lemma 2.1, (2.9) and (2.10) because $\phi(\theta) \in-\Pi_{m}$.

Let us discuss higher moments of the Wishart law $\gamma_{q, \theta}$. The computation is reduced to the derivations of the Laplace transform $L_{\mu_{q}}(\theta)=\pi^{m / 2}(\operatorname{det} \phi(-\theta))^{-1 / 2}$. Namely we have for $\eta_{1}, \ldots, \eta_{N} \in\left(\mathbb{R}^{n}\right)^{*}$

$$
E\left(\left\langle Y, \eta_{1}\right\rangle\left\langle Y, \eta_{2}\right\rangle \ldots\left\langle Y, \eta_{N}\right\rangle\right)=\frac{D_{\eta_{1}} D_{\eta_{2}} \ldots D_{\eta_{N}} L_{\mu_{q}}(\theta)}{L_{\mu_{q}}(\theta)} .
$$

On the other hand, if $f(x)=\operatorname{det}(-x)^{-p}$ for $x \in-\Pi_{m}$, where $p$ is a real constant, the derivative $D_{v_{1}} D_{v_{2}} \ldots D_{v_{N}} f(a)\left(a \in-\Pi_{m}, v_{1}, \ldots, v_{N} \in \operatorname{Sym}(m, \mathbb{R})\right)$ is given in [8, Lemma 5]. We review the formula briefly. For an element $\pi$ of the symmetric group $\mathfrak{S}_{N}$, we write $C(\pi)$ for the set of cycles of $\pi$. For $y \in \Pi_{m}$, we denote by $r_{\pi}(y)$ the multilinear form of $v_{1}, \ldots, v_{N}$ given by

$$
r_{\pi}(y)\left(v_{1}, \ldots, v_{N}\right):=\prod_{c \in C(\pi)} \operatorname{tr}\left(\prod_{j \in c} y v_{j}\right) .
$$

Then we have

$$
D_{v_{1}} D_{v_{2}} \ldots D_{v_{N}} f(a)=f(a) \cdot \sum_{\pi \in \mathfrak{G}_{N}} p^{\sharp C(\pi)} r_{\pi}\left((-a)^{-1}\right)\left(v_{1}, \ldots, v_{N}\right),
$$

which is deduced from (2.9) and (2.10) by induction (see [8]). Making use of the formula, we obtain

Theorem 2.6. Let $Y$ be an $\mathbb{R}^{n}$-valued random variable with the Wishart law $\gamma_{q, \theta}$. For $\eta_{1}, \eta_{2}, \ldots, \eta_{N} \in\left(\mathbb{R}^{n}\right)^{*}$, one has

$$
\begin{aligned}
E\left(\left\langle Y, \eta_{1}\right\rangle\left\langle Y, \eta_{2}\right\rangle \ldots\left\langle Y, \eta_{N}\right\rangle\right) & =\sum_{\pi \in \mathfrak{G}_{N}}\left(\frac{1}{2}\right)^{\sharp C(\pi)} r_{\pi}\left(\phi(-\theta)^{-1}\right)\left(\phi\left(\eta_{1}\right), \phi\left(\eta_{2}\right), \ldots, \phi\left(\eta_{N}\right)\right) \\
& =\sum_{\pi \in \mathfrak{G}_{N}}\left(\frac{1}{2}\right)^{\sharp C(\pi)} \prod_{c \in C(\pi)} \operatorname{tr}\left(\prod_{j \in c} \phi(-\theta)^{-1} \phi\left(\eta_{j}\right)\right) .
\end{aligned}
$$

2.3. Wishart laws associated to virtual quadratic maps. We shall consider virtual quadratic maps, that is a 'formal linear combination' of quadratic maps, and the associated Wishart laws. First we introduce the notion of direct sum of quadratic maps. Let $q_{i}: \mathbb{R}^{m_{i}} \rightarrow \mathbb{R}^{n}(i=1, \ldots, s)$ be $\Omega$-positive quadratic maps. Then the direct sum $q=q_{1} \oplus q_{2} \oplus \cdots \oplus q_{s}$ is an $\mathbb{R}^{n}$-valued quadratic map on
$\mathbb{R}^{m_{1}} \oplus \mathbb{R}^{m_{2}} \oplus \cdots \oplus \mathbb{R}^{m_{s}}$ given by

$$
q(x):=q_{1}\left(x_{1}\right)+q_{2}\left(x_{2}\right)+\cdots+q_{s}\left(x_{s}\right) \quad\left(x=\sum_{k=1}^{s} x_{i}, \quad x_{i} \in \mathbb{R}^{m_{i}}\right) .
$$

It is easy to see that $q$ is also $\Omega$-positive. If $q_{1}=q_{2}=\cdots=q_{s}$, then the direct sum $q$ is denoted by $q_{1}^{\oplus s}$.

The linear map $\phi:\left(\mathbb{R}^{n}\right)^{*} \rightarrow \operatorname{Sym}(m, \mathbb{R})\left(m:=\sum_{i=1}^{s} m_{i}\right)$ associated to the direct $\operatorname{sum} q=\sum^{\oplus} q_{i}$ is given by

$$
\phi(\eta)=\left(\begin{array}{cccc}
\phi_{1}(\eta) & & &  \tag{2.11}\\
& \phi_{2}(\eta) & & \\
& & \ddots & \\
& & & \phi_{s}(\eta)
\end{array}\right) \quad\left(\eta \in\left(\mathbb{R}^{n}\right)^{*}\right)
$$

Conversely, if a symmetric matrix $\phi(\eta)$ is expressed by $\phi_{1}(\eta), \ldots, \phi_{s}(\eta)$ as above for all $\eta \in\left(\mathbb{R}^{n}\right)^{*}$, then the corresponding quadratic map $q$ is the direct sum of $q_{1}, \ldots, q_{s}$. In Example 1, the quadratic map $q: \mathbb{R}^{4} \rightarrow \mathbb{R}^{3}$ is the direct sum of 4 quadratic maps $q_{i}: \mathbb{R} \ni x \mapsto x^{2} v_{i} \in \mathbb{R}^{3}(i=1, \ldots, 4)$, where

$$
v_{1}:=\left(\begin{array}{l}
0 \\
0 \\
1
\end{array}\right), \quad v_{2}:=\left(\begin{array}{l}
1 \\
0 \\
1
\end{array}\right), \quad v_{3}:=\left(\begin{array}{l}
1 \\
1 \\
1
\end{array}\right), \quad q_{4}:=\left(\begin{array}{l}
0 \\
1 \\
1
\end{array}\right) .
$$

In Example 2, we see that $q_{r, s}: \operatorname{Mat}(r, s ; \mathbb{R}) \rightarrow \operatorname{Sym}(r, \mathbb{R})$ is naturally identified with $q_{r, 1}^{\oplus s}$. In Example 3 with $\mathcal{Z}$ given by (2.4), we have $q_{\mathcal{Z}}=q_{\mathcal{Z}}^{\{1\}} \oplus q_{\mathcal{Z}}^{\{2,3\}}$, while we do not have such a decomposition for the case (2.5).

Let $q_{i}: \mathbb{R}^{m_{i}} \rightarrow \mathbb{R}^{n}(i=1,2)$ be $\Omega$-positive quadratic maps, and $q$ the direct sum $q_{1} \oplus q_{2}$. Then it is easy to see that the measure $\mu_{q}$ equals the convolution $\mu_{q_{1}} * \mu_{q_{2}}$. Thus, for $\theta \in-\Omega^{*}$ we have $L_{\mu_{q}}(\theta)=L_{\mu_{q_{1}}}(\theta) L_{\mu_{q_{2}}}(\theta)$ and $\gamma_{q_{,}, \theta}=\gamma_{q_{1}, \theta} * \gamma_{q_{2}, \theta}$. In general, if we set $q=q_{1}^{\oplus s_{1}} \oplus q_{2}^{\oplus s_{2}} \oplus \cdots \oplus q_{t}^{\oplus s_{t}}$ for $\Omega$-positive quadratic maps $q_{i}: \mathbb{R}^{m_{i}} \rightarrow \mathbb{R}^{n}(i=1,2, \ldots, t)$ and positive integers $s_{1}, s_{2}, \ldots, s_{t}$, then we have

$$
\begin{equation*}
L_{\mu_{q}}(\theta)=\prod_{i=1}^{t} L_{\mu_{i}}(\theta)^{s_{i}} \quad\left(\theta \in-\Omega^{*}\right) \tag{2.12}
\end{equation*}
$$

Now we remark that, even though $s_{i}$ 's are not positive integers, there may exist a positive measure $\mu_{q}$ on $\bar{\Omega}$ for which the relation (2.12) holds. In general, for real numbers $s_{1}, \ldots, s_{r}$, we shall call a formal sum $q=q_{1}^{\oplus s_{1}} \oplus q_{2}^{\oplus s_{2}} \oplus \cdots \oplus q_{t}^{\oplus s_{t}}$ a virtual $\Omega$ positive quadratic map. The measure $\mu_{q}$ satisfying (2.12) is called the Riesz measure associated to $q$. By the injectivity of the Laplace transform, the associated Riesz measure $\mu_{q}$ is unique if it exists. In this case, the Wishart laws $\gamma_{q, \theta}\left(\theta \in-\Omega^{*}\right)$ are defined again as members of the exponential family generated by $\mu_{q}$.

Proposition 2.7. Let $q_{i}: \mathbb{R}^{m_{i}} \rightarrow \mathbb{R}^{n}(i=1, \ldots, t)$ be $\Omega$-positive quadratic maps. Assume that there exists a measure $\mu_{q}$ associated to the virtual quadratic map $q=$ $q_{1}^{\oplus s_{1}} \oplus \cdots \oplus q_{t}^{\oplus s_{t}}$ for certain $s_{1}, \ldots, s_{t} \in \mathbb{R}$. Let $Y$ be an $\mathbb{R}^{n}$-valued random variable with the Wishart law $\gamma_{q, \theta}$. Then the Laplace transform $L_{\gamma_{q, \theta}}(\eta)=E\left(e^{\langle Y, \eta\rangle}\right)$ of the law $\gamma_{q, \theta}$ is given by

$$
L_{\gamma_{q, \theta}}(\eta)=\prod_{i=1}^{t} L_{\gamma_{q_{i}, \theta}}(\eta)^{s_{i}}=\prod_{i=1}^{t} \operatorname{det}\left(I_{m_{i}}+\phi_{i}(-\theta)^{-1} \phi_{i}(-\eta)\right)^{-s_{i} / 2}
$$

for $\eta \in-\theta-\Omega^{*}$.

Proof. We have $L_{\gamma_{q, \theta}}(\eta)=L_{\mu_{q}}(\theta)^{-1} L_{\mu_{q}}(\eta+\theta)$ by definition, and the right-hand side equals $\prod_{i=1}^{t}\left(L_{\mu_{i}}(\theta)^{-1} L_{\mu_{i}}(\eta+\theta)\right)^{s_{i}}$ by (2.12). Since $L_{\gamma_{q_{i}, \theta}}(\eta)=L_{\mu_{i}}(\theta)^{-1} L_{\mu_{i}}(\eta+$ $\theta$ ), we obtain the first equality. The second equality follows from Proposition 2.4.

Since we see immediately from (2.12) that

$$
\log L_{\mu_{q}}(\theta)=\sum_{i=1}^{t} s_{i} \log L_{\mu_{i}}(\theta)
$$

the virtual versions of Theorem 2.5 is given as follows:

Proposition 2.8. Under the same assumption of Proposition 2.7, one has
(i) $E(\langle Y, \eta\rangle)=\sum_{i=1}^{t} s_{i} \operatorname{tr} \phi_{i}(-\theta)^{-1} \phi_{i}(\eta) / 2$ for $\eta \in\left(\mathbb{R}^{n}\right)^{*}$,
(ii) $E\left((\langle Y, \eta\rangle-M)\left(\left\langle Y, \eta^{\prime}\right\rangle-M^{\prime}\right)\right)=\sum_{i=1}^{t} s_{i} \operatorname{tr} \phi_{i}(-\theta)^{-1} \phi_{i}(\eta) \phi_{i}(-\theta)^{-1} \phi_{i}\left(\eta^{\prime}\right) / 2$ for $\eta, \eta^{\prime} \in\left(\mathbb{R}^{n}\right)^{*}$, where $M:=E(\langle Y, \eta\rangle)$ and $M^{\prime}:=E\left(\left\langle Y, \eta^{\prime}\right\rangle\right)$.

As for higher moments, we generalize the formula in Theorem 2.6 as follows:
Theorem 2.9. Under the same assumption of Proposition 2.7, one has

$$
\begin{align*}
& E\left(\left\langle Y, \eta_{1}\right\rangle\left\langle Y, \eta_{2}\right\rangle \ldots\left\langle Y, \eta_{N}\right\rangle\right) \\
& =\sum_{\pi \in \mathfrak{S}_{N}}\left(\frac{1}{2}\right)^{\sharp C(\pi)} \prod_{c \in C(\pi)}\left\{\sum_{i=1}^{t} s_{i} \operatorname{tr}\left(\prod_{j \in c} \phi_{i}(-\theta)^{-1} \phi_{i}\left(\eta_{j}\right)\right)\right\} . \tag{2.13}
\end{align*}
$$

for $\eta_{1}, \eta_{2}, \ldots, \eta_{N} \in\left(\mathbb{R}^{n}\right)^{*}$.
Theorem 2.9 easily follows from Theorem 2.6 when $s_{1}, \ldots, s_{t}$ are positive integers, that is, $q$ is a true quadratic map. Indeed, we can consider the associated linear
$\operatorname{map} \phi=\phi_{q}$ in this case, and we have

$$
\operatorname{tr}\left(\prod_{j \in c} \phi(-\theta)^{-1} \phi\left(\eta_{j}\right)\right)=\sum_{i=1}^{t} s_{i} \operatorname{tr}\left(\prod_{j \in c} \phi_{i}(-\theta)^{-1} \phi_{i}\left(\eta_{j}\right)\right)
$$

by virtue of (2.11). To prove (2.13) for general case, it is enough to verify that the quantity $E\left(\left\langle Y, \eta_{1}\right\rangle\left\langle Y, \eta_{2}\right\rangle \ldots\left\langle Y, \eta_{N}\right\rangle\right)$ is a polynomial of $s_{1}, \ldots, s_{t}$. For this purpose, we make some calculations involving the semi-invariants or the cummulants (cf. [18]).

For $\eta \in\left(\mathbb{R}^{n}\right)^{*}, \theta \in-\Omega^{*}$ and a positive integer $k$, we denote by $S_{k}(\theta ; \eta)$ the $k$-th derivative of $\log L_{\mu_{q}}(\theta)$ in the direction $\eta$ :

$$
S_{k}(\theta ; \eta)=\left(D_{\eta}\right)^{k} \log L_{\mu_{q}}(\theta)=\left(\frac{d^{k}}{d h^{k}}\right)_{h=0} \log L_{\mu_{q}}(\theta+h \eta)
$$

Similarly to Proposition 2.8, we get

$$
\begin{equation*}
S_{k}(\theta ; \eta)=\frac{(k-1)!}{2} \sum_{i=1}^{t} s_{i} \operatorname{tr}\left(\phi_{i}(-\theta)^{-1} \phi_{i}(\eta)\right)^{k} \tag{2.14}
\end{equation*}
$$

On the other hand, since the function $L_{\mu_{q}}(\theta)$ is analytic, we have

$$
\log L_{\mu_{q}}(\theta+h \eta)=\log L_{\mu_{q}}(\theta)+\sum_{k=1}^{\infty} \frac{h^{k}}{k!} S_{k}(\theta ; \eta)
$$

Taking the exponential, we have

$$
\begin{align*}
L_{\mu_{q}}(\theta+h \eta) & =L_{\mu_{q}}(\theta) \exp \left(\sum_{k=1}^{\infty} \frac{h^{k}}{k!} S_{k}(\theta ; \eta)\right)  \tag{2.15}\\
& =L_{\mu_{q}}(\theta) \sum_{\ell=0}^{\infty} \frac{1}{\ell!}\left(\sum_{k=1}^{\infty} \frac{h^{k}}{k!} S_{k}(\theta ; \eta)\right)^{\ell} .
\end{align*}
$$

Let $Y$ be an $\mathbb{R}^{n}$-valued random variable with the Wishart distribution $\gamma_{q, \theta}$. Then we have

$$
E\left(\langle Y, \eta\rangle^{N}\right)=\frac{\left(D_{\eta}\right)^{N} L_{\mu_{q}}(\theta)}{L_{\mu_{q}}(\theta)}
$$

so that

$$
\begin{equation*}
L_{\mu_{q}}(\theta+h \eta)=L_{\mu_{q}}(\theta) \sum_{N=0}^{\infty} \frac{h^{N}}{N!} E\left(\langle Y, \eta\rangle^{N}\right) \tag{2.16}
\end{equation*}
$$

Comparing the coefficients of $h^{N}$ in (2.15) and (2.16), we obtain

$$
E\left(\langle Y, \eta\rangle^{N}\right)=\sum_{\ell=1}^{N} \frac{1}{\ell!} \sum_{k_{1}+k_{2}+\cdots+k_{\ell}=N} \frac{N!}{k_{1}!k_{2}!\cdots k_{\ell}!} S_{k_{1}}(\theta ; \eta) S_{k_{2}}(\theta ; \eta) \ldots S_{k_{\ell}}(\theta ; \eta)
$$

Substituting (2.14) to the equality above, we obtain a formula for the higher moments.

Proposition 2.10. One has

$$
\begin{equation*}
E\left(\langle Y, \eta\rangle^{N}\right)=\sum_{\ell=1}^{N} \frac{1}{\ell!} \sum_{k_{1}+k_{2}+\cdots+k_{\ell}=N} \frac{N!}{k_{1} k_{2} \cdots k_{\ell}} \prod_{j=1}^{\ell}\left(\sum_{i=1}^{t} \frac{s_{i}}{2} \operatorname{tr}\left(\phi_{i}(-\theta)^{-1} \phi_{i}(\eta)\right)^{k_{j}}\right) . \tag{2.17}
\end{equation*}
$$

The moments $E\left(\left\langle Y, \eta_{1}\right\rangle\left\langle Y, \eta_{2}\right\rangle \cdots\left\langle Y, \eta_{N}\right\rangle\right)\left(\eta_{1}, \ldots, \eta_{N} \in\left(\mathbb{R}^{n}\right)^{*}\right)$ can be obtained now by polarization, so that it is a polynomial of $s_{1}, \ldots, s_{t}$. Hence the proof of Theorem 2.9 is completed. We note that (2.13) with $\eta_{1}=\cdots=\eta_{N}=\eta$ becomes (2.17), which means that Proposition 2.10 together with the polarization process also yields Theorem 2.9.
2.4. Group equivariance of the Wishart laws. Let $G(\Omega)$ be the linear automorphism group $\{g \in G L(n, \mathbb{R}) ; g \Omega=\Omega\}$ of $\Omega$. For an $\Omega$-positive quadratic map $q: \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$ and $g \in G(\Omega)$, the quadratic map $g \circ q: \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$ is again $\Omega$-positive. It is easy to see that the Riesz measure $\mu_{g \circ q}$ is the image of $\mu_{q}$ by $g$, that is,

$$
\begin{equation*}
\mu_{g \circ q}(A)=\mu_{q}\left(g^{-1} A\right) \tag{2.18}
\end{equation*}
$$

for a measurable set $A \subset \mathbb{R}^{n}$. Let us discuss the Wishart laws $\gamma_{g \circ q, \theta}$ for $\theta \in-\Omega^{*}$. For $\eta \in\left(\mathbb{R}^{n}\right)^{*}$, we denote by $g^{*} \eta$ the linear form $\eta \circ g \in\left(\mathbb{R}^{n}\right)^{*}$. If $\eta \in \Omega^{*}$, then $g^{*} \eta \in \Omega^{*}$ because $\left\langle y, g^{*} \eta\right\rangle=\langle g y, \eta\rangle>0$ for $y \in \bar{\Omega} \backslash\{0\}$. We observe

$$
\begin{equation*}
L_{\mu_{q}}\left(g^{*} \theta\right)=\int_{\mathbb{R}^{m}} e^{\left\langle q(x), g^{*} \theta\right\rangle} d x=\int_{\mathbb{R}^{m}} e^{\langle g \circ q(x), \theta\rangle} d x=L_{\mu_{g \circ q}}(\theta) . \tag{2.19}
\end{equation*}
$$

Therefore, denoting by $1_{A}$ the characteristic function of a measurable set $A \subset \mathbb{R}^{n}$, we have

$$
\begin{align*}
\gamma_{g \circ q, \theta}(A) & =\frac{1}{L_{\mu_{g \circ q}}(\theta)} \int_{\mathbb{R}^{m}} 1_{A}(g \circ q(x)) e^{\langle g \circ q(x), \theta\rangle} d x  \tag{2.20}\\
& =\frac{1}{L_{\mu_{q}}\left(g^{*} \theta\right)} \int_{\mathbb{R}^{m}} 1_{g^{-1} A}(q(x)) e^{\left\langle q(x), g^{*} \theta\right\rangle} d x=\gamma_{q, g^{*} \theta}\left(g^{-1} A\right) .
\end{align*}
$$

We restate (2.20) as follows.
Lemma 2.11. Let $g$ be an element of $G(\Omega)$. If a random variable $Y$ obeys the Wishart law $\gamma_{q, \theta}$, the law of $g Y$ is $\gamma_{g \circ q,\left(g^{-1}\right)^{*} \theta}$.

Let $q_{i}: \mathbb{R}^{m_{i}} \rightarrow \mathbb{R}^{n}(i=1, \ldots, t)$ be $\Omega$-positive quadratic maps, and $q$ the virtual quadratic map $q_{1}^{\oplus s_{1}} \oplus \cdots \oplus q_{t}^{\oplus s t}$ with $s_{1}, \ldots, s_{t} \in \mathbb{R}$. Then we define $g \circ q$ to be the virtual quadratic map $\left(g \circ q_{1}\right)^{\oplus s_{1}} \oplus \cdots \oplus\left(g \circ q_{t}\right)^{\oplus s_{t}}$.

Proposition 2.12. If the Riesz measure $\mu_{q}$ exists, then the Riesz measure $\mu_{g \circ q}$ exists and equals the image of $\mu_{q}$ by $g$. Moreover $\gamma_{g \circ q,\left(g^{-1}\right)^{* \theta}}$ is the image of $\gamma_{q, \theta}$ by $g$.

Proof. Let $\mu^{\prime}$ be the image of $\mu_{q}$ by $g$. For $\theta \in-\Omega^{*}$, we have

$$
\begin{equation*}
L_{\mu^{\prime}}(\theta)=\int_{\mathbb{R}^{n}} e^{\langle y, \theta\rangle} \mu^{\prime}(d y)=\int_{\mathbb{R}^{n}} e^{\langle g y, \theta\rangle} \mu_{q}(d y)=\int_{\mathbb{R}^{n}} e^{\left\langle y, g^{*} \theta\right\rangle} \mu_{q}(d y)=L_{\mu_{q}}\left(g^{*} \theta\right) \tag{2.21}
\end{equation*}
$$

By (2.12) and (2.19), the last term equals $\prod_{i=1}^{t} L_{\mu_{i}}\left(g^{*} \theta\right)^{s_{i}}=\prod_{i=1}^{t} L_{\mu_{g \circ q_{i}}}(\theta)^{s_{i}}$. Thus we get $L_{\mu^{\prime}}(\theta)=\prod_{i=1}^{t} L_{\mu_{g \circ q_{i}}}(\theta)^{s_{i}}$ which means $\mu^{\prime}=\mu_{g \circ q}$ by (2.12).

Let $\gamma^{\prime}$ be the image of $\gamma_{q, \theta}$ by $g$. Similarly to (2.21), we have $L_{\gamma^{\prime}}(\eta)=L_{\gamma_{q, \theta}}\left(g^{*} \eta\right)$ for $\eta \in-\left(g^{-1}\right)^{*} \theta-\Omega^{*}$, while $L_{\gamma_{g \circ q_{i},\left(g^{-1}\right)^{*} \theta}}(\eta)=L_{\gamma_{q_{i}, \theta}}\left(g^{*} \eta\right)$ by Lemma 2.11. On the other hand, we see from Proposition 2.7 that

$$
L_{\gamma_{q, \theta}}\left(g^{*} \eta\right)=\prod_{i=1}^{t} L_{\gamma_{q_{i}, \theta}}\left(g^{*} \eta\right)^{s_{i}}=\prod_{i=1}^{t} L_{\gamma_{g \circ q_{i},\left(g^{-1}\right)^{* \theta}}}(\eta)^{s_{i}}=L_{\gamma_{g \circ q,\left(g^{-1}\right)^{* \theta}}}(\eta)
$$

Thus we get $L_{\gamma^{\prime}}(\eta)=L_{\gamma_{g \circ q,\left(g^{-1}\right)^{*} \theta}}(\eta)$, so that $\gamma^{\prime}=\gamma_{g \circ q,\left(g^{-1}\right)^{* \theta}}$ by the injectivity of the Laplace transform.

Let $\operatorname{Aut}(\Omega, q)$ be the set of pairs $\left(g_{1}, g_{2}\right) \in G(\Omega) \times G L(m, \mathbb{R})$ for which $g_{1} \circ q=q \circ g_{2}$. Then $\operatorname{Aut}(\Omega, q)$ forms a Lie subgroup of $G L(n, \mathbb{R}) \times G L(m, \mathbb{R})$, and we have a group homomorphism

$$
p r_{1}: \operatorname{Aut}(\Omega, q) \ni\left(g_{1}, g_{2}\right) \mapsto g_{1} \in G(\Omega)
$$

The condition $\left(g_{1}, g_{2}\right) \in \operatorname{Aut}(\Omega, q)$ is also equivalent to

$$
\begin{equation*}
\phi_{q}\left(g_{1}^{*} \eta\right)={ }^{\mathrm{t}} g_{2} \phi_{q}(\eta) g_{2} \quad\left(\eta \in\left(\mathbb{R}^{n}\right)^{*}\right) \tag{2.22}
\end{equation*}
$$

Then we obtain

$$
\begin{equation*}
\operatorname{det} \phi_{q}\left(g_{1}^{*} \eta\right)=C \operatorname{det} \phi_{q}(\eta) \quad\left(\eta \in\left(\mathbb{R}^{n}\right)^{*}\right) \tag{2.23}
\end{equation*}
$$

with $C=\left(\operatorname{det} g_{2}\right)^{2}$, which means that $\operatorname{det} \phi_{q}(\eta)$ is a relatively invariant polynomial on $\left(\mathbb{R}^{n}\right)^{*}$ under the contragredient action of $\operatorname{pr}_{1}(\operatorname{Aut}(\Omega, q))$. The following proposition describes a transformation rule of the family of the Wishart laws $\left\{\gamma_{q, \theta}\right\}_{\theta \in-\Omega^{*}}$ under the group $\operatorname{pr}_{1}(\operatorname{Aut}(\Omega, q))$.

Proposition 2.13. For a measurable set $A \subset \mathbb{R}^{n}$ and $\left(g_{1}, g_{2}\right) \in \operatorname{Aut}(\Omega, q)$, one has (i) $\mu_{q}\left(g_{1}^{-1} A\right)=\mu_{g_{1} \circ q}(A)=\left|\operatorname{det} g_{2}\right|^{-1} \mu_{q}(A)$, (ii) $\gamma_{q, g_{1}^{*} \theta}(A)=\gamma_{q, \theta}\left(g_{1} A\right)$.

Proof. (i) Because of (2.18), we only have to show the second equality. By definition, we have

$$
\mu_{g_{1} \circ q}(A)=\int_{\mathbb{R}^{m}} 1_{A}\left(g_{1} \circ q(x)\right) d x=\int_{\mathbb{R}^{m}} 1_{A}\left(q \circ g_{2}(x)\right) d x
$$

Putting $x^{\prime}=g_{2} x$, the last term equals

$$
\left|\operatorname{det} g_{2}\right|^{-1} \int_{\mathbb{R}^{m}} 1_{A}\left(q\left(x^{\prime}\right)\right) d x^{\prime}=\left|\operatorname{det} g_{2}\right|^{-1} \mu_{q}(A)
$$

whence (i) follows.
(ii) By (2.20), we get for $y \in \mathbb{R}^{n}$

$$
\gamma_{q, g_{1}^{*} \theta}(d y)=\gamma_{g_{1} \circ q, \theta}\left(g_{1} d y\right)=\frac{e^{\left\langle g_{1} y, \theta\right\rangle}}{L_{\mu_{g_{1} \circ q}}(\theta)} \mu_{g_{1} \circ q}\left(g_{1} d y\right),
$$

Since $\mu_{g_{1} \circ q}=\left|\operatorname{det} g_{2}\right|^{-1} \mu_{q}$ by (i), the last term equals

$$
\frac{e^{\left\langle g_{1} y, \theta\right\rangle}}{\left|\operatorname{det} g_{2}\right|^{-1} L_{\mu_{q}}(\theta)}\left|\operatorname{det} g_{2}\right|^{-1} \mu_{q}\left(g_{1} d y\right)=\gamma_{q, \theta}\left(g_{1} d y\right)
$$

Hence (ii) is verified.

## §3. Homogeneous Case.

3.1. Homogeneous quadratic map. An $\Omega$-positive map $q: \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$ is said to be homogeneous if, for any $y, y^{\prime} \in \Omega$, there exists $\left(g_{1}, g_{2}\right) \in \operatorname{Aut}(\Omega, q)$ for which $g_{1} y=y^{\prime}$. In other words, $q$ is homogeneous if $\operatorname{pr}_{1}(\operatorname{Aut}(q, \Omega))$ acts on $\Omega$ transitively. In this case, $\Omega$ is clearly a homogeneous cone, that is, a linear group on $\mathbb{R}^{n}$ acts on the cone $\Omega$ transitively. Then the dual cone $\Omega^{*} \subset\left(\mathbb{R}^{n}\right)^{*}$ is also a homogeneous cone on which the group $\operatorname{pr}_{1}(\operatorname{Aut}(q, \Omega))$ acts transitively by the contragredient action ([23]). We see from (2.1) and (2.22) that the quadratic map $q$ is homogeneous if and only if the associated linear map $\phi_{q}:\left(\mathbb{R}^{n}\right)^{*} \rightarrow \operatorname{Sym}(m, \mathbb{R})$ is a representation of the dual cone $\Omega^{*}$ in the sense of Rothaus [22] (see also [15]).

A typical example of a homogeneous cone is $\Pi_{r} \subset \operatorname{Sym}(r, \mathbb{R})$. For $A \in G L(r, \mathbb{R})$, we denote by $\rho(A)$ the linear map on $\operatorname{Sym}(r, \mathbb{R})$ defined by $\rho(A) y:=A y^{\mathrm{t}} A(y \in$ $\operatorname{Sym}(r, \mathbb{R}))$. Then the group $\rho(G L(r, \mathbb{R}))$ acts on $\Pi_{r}$ transitively. Moreover, the linear automorphism group $G\left(\Pi_{r}\right)$ equals $\rho(G L(r, \mathbb{R}))$. We see that the quadratic map $q_{r, s}: \operatorname{Mat}(r, s ; \mathbb{R}) \rightarrow \operatorname{Sym}(r, \mathbb{R})$ in Example 2 is homogeneous. Indeed, we have a surjective homomorphism $G L(r, \mathbb{R}) \times O(s) \ni(A, B) \mapsto\left(\rho(A), \tau_{r, s}(A, B)\right) \in$ $\operatorname{Aut}\left(\Pi_{r}, q_{r, s}\right)$, where $\tau_{r, s}(A, B)$ is a linear map on $\operatorname{Mat}(r, s ; \mathbb{R})$ given by $\tau_{r, s}(A, B) x:=$
$\operatorname{Ax} B^{-1}(x \in \operatorname{Mat}(r, s ; \mathbb{R}))$. Thus we have $\operatorname{pr}_{1}\left(\operatorname{Aut}\left(q_{r, s}, \Pi_{r}\right)\right)=\rho(G L(r, \mathbb{R}))$, which acts on $\Pi_{r}$ transitively.

For a subset $I \subset\{1, \ldots, r\}$, we denote by $q^{I}$ the restriction of $q_{r, 1}: \mathbb{R}^{r} \rightarrow$ $\operatorname{Sym}(r, \mathbb{R})$ to the space $R^{I} \subset \mathbb{R}^{r}$ defined in (2.6). The map $q^{I}$ coincides with $q_{\mathcal{Z}}^{I}$ in Example 3 with $\mathcal{Z}=\operatorname{Sym}(r, \mathbb{R})$. Let us observe that $q^{I}$ is homogeneous in general. Let $P^{I}$ be the linear group consisting of $A \in G L(r, \mathbb{R})$ for which $A R^{I}=R^{I}$. For example, if $r=3$, we have

$$
\begin{align*}
P^{\{1\}} & =\left\{A=\left(\begin{array}{ccc}
a_{11} & a_{12} & a_{13} \\
0 & a_{22} & a_{23} \\
0 & a_{32} & a_{33}
\end{array}\right) ; A \in G L(3, \mathbb{R})\right\},  \tag{3.1}\\
P^{\{2,3\}} & =\left\{A=\left(\begin{array}{ccc}
a_{11} & 0 & 0 \\
a_{21} & a_{22} & a_{23} \\
a_{31} & a_{32} & a_{33}
\end{array}\right) ; A \in G L(3, \mathbb{R})\right\} . \tag{3.2}
\end{align*}
$$

Since we have a homomorphism $P^{I} \ni A \mapsto(\rho(A), A) \in \operatorname{Aut}\left(\Pi_{r}, q^{I}\right)$, it is enough to show that $\rho\left(P^{I}\right)$ acts on $\Pi_{r}$ transitively. Put $k:=\sharp I$ and take a permutation matrix $w_{0} \in \mathfrak{S}_{r} \subset G L(r, \mathbb{R})$ sending $R^{\{r-k+1, \ldots, r\}}$ onto $R^{I}$. Then we have $P^{I}=$ $w_{0} P^{\{r-k+1, \ldots, r\}} w_{0}^{-1}$, and

$$
P^{\{r-k+1, \ldots, r\}}=\left\{\left(\begin{array}{cc}
A_{1} & 0 \\
A_{2} & A_{3}
\end{array}\right) ; \begin{array}{l}
A_{1} \in G L(k, \mathbb{R}), \quad A_{2} \in \operatorname{Mat}(r-k+1, k ; \mathbb{R}) \\
A_{3} \in G L(r-k+1, \mathbb{R})
\end{array}\right\} .
$$

Since $P^{\{r-k+1, \ldots, r\}}$ contains the group of lower triangular matrices, $\rho\left(P^{\{r-k+1, \ldots, r\}}\right)$ acts on $\Pi_{r}$ transitively. Therefore $\rho\left(P^{I}\right)=\rho\left(w_{0}\right) \rho\left(P^{\{r-k+1, \ldots, r\}}\right) \rho\left(w_{0}\right)^{-1}$ also acts on $\Pi_{r}$ transitively, so that $q^{I}$ is homogeneous.

Coming back to the examples (3.1) and (3.2), we note that $q=q^{\{1\}} \oplus q^{\{2,3\}}$ is not homogeneous as $\Pi_{3}$-positive quadratic map, while both $q^{\{1\}}$ and $q^{\{2,3\}}$ are. Indeed, the image of $q$ generates the space $\mathcal{Z} \subset \operatorname{Sym}(3, \mathbb{R})$ in (2.4). Thus, if $\left(g_{1}, g_{2}\right) \in$ $\operatorname{Aut}\left(\Pi_{3}, q\right)$, then $g_{1}$ must preserve both $\mathcal{Z}$ and $\Pi_{3}$. Let us take $y \in \Pi_{3} \backslash \mathcal{Z}$. Then $g_{1}$ does not send $I_{3} \in \Pi_{3}$ to $y$ because $I_{3} \in \mathcal{Z}$. Thus the action of $\operatorname{pr}_{1}\left(\operatorname{Aut}\left(q, \Pi_{3}\right)\right)$ on $\Pi_{3}$ is not transitive.

On the other hand, if we regard $q$ as a map from $R^{\{1\}} \oplus R^{\{2,3\}}$ to $\mathcal{Z}$, then $q$ is a homogeneous $\mathcal{P}$-positive quadratic map, where $\mathcal{P}:=\mathcal{Z} \cap \Pi_{r}$. In fact, since $(\rho(A), A) \in \operatorname{Aut}\left(q^{\{1\}}, \mathcal{P}\right) \cap \operatorname{Aut}\left(q^{\{2,3\}}, \mathcal{P}\right)$ for $A \in P^{\{1\}} \cap P^{\{2,3\}}$, we have $\rho(A) \in$ $p r_{1}(\operatorname{Aut}(q, \mathcal{P}))$. Therefore $\operatorname{pr} r_{1}(\operatorname{Aut}(q, \mathcal{P}))$ contains a group $\rho\left(P^{\{1\}} \cap P^{\{2,3\}}\right)$ which acts on $\mathcal{P}$ transitively.

In Example 1, the quadratic map $q: \mathbb{R}^{4} \rightarrow \mathbb{R}^{3}$ is not homogeneous because $\Omega \subset \mathbb{R}^{4}$ in (2.2) is not a homogeneous cone ([13]).
3.2. Matrix realization of homogeneous cones. In this section, we shall discuss a homogeneous cone realized as $\mathcal{P}_{\mathcal{V}}=\mathcal{Z}_{\mathcal{V}} \cap \Pi_{N}$ with $\mathcal{Z}_{\mathcal{V}} \subset \operatorname{Sym}(N, \mathbb{R})$ constructed from an appropriate system $\mathcal{V}=\left\{\mathcal{V}_{l k}\right\}$ of vector spaces in a specific way explained below, following [14, section 3.1]. The investigation of such cones is fundamental because all homogeneous cones are linearly equivalent to some $\mathcal{P}_{\mathcal{V}}$ due to [14, Theorem D].

Let us take a partition $N=n_{1}+n_{2}+\cdots+n_{r}$ of a positive integer $N$, and consider a system of vector spaces $\mathcal{V}_{l k} \subset \operatorname{Mat}\left(n_{l}, n_{k} ; \mathbb{R}\right)(1 \leq k<l \leq r)$ satisfying the following three conditions:
(V1) $A \in \mathcal{V}_{l k}, B \in \mathcal{V}_{k j} \Rightarrow A B \in \mathcal{V}_{l j} \quad(1 \leq j<k<l \leq r)$,
(V2) $A \in \mathcal{V}_{l j}, B \in \mathcal{V}_{k j} \Rightarrow A^{\mathrm{t}} B \in \mathcal{V}_{l k} \quad(1 \leq j<k<l \leq r)$,
(V3) $A \in \mathcal{V}_{l k} \Rightarrow A^{\mathrm{t}} A \in \mathbb{R} I_{n_{l}} \quad(1 \leq k<l \leq r)$.
Let $\mathcal{Z}_{\mathcal{V}}$ be the subspace of $\operatorname{Sym}(N, \mathbb{R})$ defined by

$$
\mathcal{Z}_{\mathcal{V}}:=\left\{y=\left(\begin{array}{cccc}
Y_{11} & { }^{\mathrm{t}} Y_{21} & \cdots & { }^{\mathrm{t}} Y_{r 1} \\
Y_{21} & Y_{22} & & { }^{\mathrm{t}} Y_{r 2} \\
\vdots & & \ddots & \\
Y_{r 1} & Y_{r 2} & \cdots & Y_{r r}
\end{array}\right) ; \begin{array}{l}
Y_{k k}=y_{k k} I_{n_{k}}, y_{k k} \in \mathbb{R}(k=1, \ldots, r) \\
Y_{l k} \in \mathcal{V}_{l k}(1 \leq k<l \leq r)
\end{array}\right\} .
$$

We set $\mathcal{P}_{\mathcal{V}}:=\mathcal{Z}_{\mathcal{V}} \cap \Pi_{N}$. Then $\mathcal{P}_{\mathcal{V}}$ is a regular open convex cone in the vector space $\mathcal{Z}_{\mathcal{V}}$. Let $H_{N}$ be the group of real lower triangular matrices with positive diagonals, and $H_{\mathcal{V}}$ a Lie subgroup of $H_{N}$ defined by

$$
H_{\mathcal{V}}:=\left\{T=\left(\begin{array}{cccc}
T_{11} & & & \\
T_{21} & T_{22} & & \\
\vdots & & \ddots & \\
T_{r 1} & T_{r 2} & \cdots & T_{r r}
\end{array}\right) ; \begin{array}{l}
T_{k k}=t_{k k} I_{n_{k}}, t_{k k}>0(k=1, \ldots, r) \\
T_{l k} \in \mathcal{V}_{l k}(1 \leq k<l \leq r)
\end{array}\right\}
$$

If $T \in H_{\mathcal{V}}$ and $y \in \mathcal{Z}_{\mathcal{V}}$, then $\rho(T) y=T y^{\mathrm{t}} T$ belongs to $\mathcal{Z}_{\mathcal{V}}$ thanks to (V1)-(V3). Moreover $\rho\left(H_{\mathcal{V}}\right)$ acts on the cone $\mathcal{P}_{\mathcal{V}} \subset \mathcal{Z}_{\mathcal{V}}$ simply transitively (cf. [14, Proposition 2.1]).

Keeping (V3) in mind, we define an inner product on the vector space $\mathcal{V}_{l k}(1 \leq$ $k<l \leq r$ ) by the equality

$$
\begin{equation*}
A^{\mathrm{t}} A=(A \mid A) I_{n_{l}}=\|A\|^{2} I_{n_{l}} \quad\left(A \in \mathcal{V}_{l k}\right) \tag{3.3}
\end{equation*}
$$

For $y, y^{\prime} \in \mathcal{Z}_{\mathcal{V}}$, we set

$$
\begin{equation*}
\left\langle y, y^{\prime}\right\rangle:=\sum_{k=1}^{r} y_{k k} y_{k k}^{\prime}+2 \sum_{1 \leq k<l \leq r}\left(Y_{l k} \mid Y_{l k}^{\prime}\right) \tag{3.4}
\end{equation*}
$$

where $y_{k k}$ and $Y_{l k}$ (respectively $y_{k k}^{\prime}$ and $Y_{l k}^{\prime}$ ) denote the components of $y$ (respectively $\left.y^{\prime}\right)$. Note that the inner product is not equal to $\operatorname{tr} y y^{\prime}$ unless $n_{1}=\cdots=n_{r}=1$. By this coupling, we identify the dual space $\mathcal{Z}_{\mathcal{V}}^{*}$ with $\mathcal{Z}_{\mathcal{V}}$. Let us observe that $I_{N}$ belongs to the dual cone $\mathcal{P}_{\mathcal{V}}^{*}$ of $\mathcal{P}_{\mathcal{V}}$, that is,

$$
0<\left\langle y, I_{N}\right\rangle=y_{11}+\cdots+y_{r r} \quad\left(y \in \overline{\mathcal{P}}_{\mathcal{V}} \backslash\{0\}\right) .
$$

Indeed, since each $y_{k k}$ is a diagonal entry of the non-negative matrix $y \in \overline{\mathcal{P}}_{\mathcal{V}} \backslash\{0\}$, we have $y_{k k} \geq 0$. Suppose $\sum_{k=1}^{r} y_{k k}=0$. Then $y_{11}=\cdots=y_{r r}=0$ and $\operatorname{tr} y=$ $\sum_{k=1}^{r} n_{k} y_{k k}=0$. This together with the non-negativity of $y$ implies $y=0$, which is a contradiction.

For $T \in H_{\mathcal{V}}$, define $\rho^{*}(T) \in G L\left(\mathcal{Z}_{\mathcal{V}}\right)$ by $\left\langle y, \rho^{*}(T) \eta\right\rangle=\langle\rho(T) y, \eta\rangle\left(y, \eta \in \mathcal{Z}_{\mathcal{V}}\right)$. By [24], we have $\mathcal{P}_{\mathcal{V}}^{*}=\left\{\rho^{*}(T) I_{N} ; T \in H_{\mathcal{V}}\right\}$. Moreover, the map $H_{\mathcal{V}} \ni T \mapsto \rho^{*}(T) I_{N} \in$ $\mathcal{P}_{\mathcal{V}}^{*}$ is a diffeomorphism. For $\underline{\sigma}=\left(\sigma_{1}, \ldots, \sigma_{r}\right) \in \mathbb{C}^{r}$, we define the one-dimensional representation $\chi_{\underline{\sigma}}: H_{\mathcal{V}} \rightarrow \mathbb{C}^{\times}$by $\chi_{\underline{\sigma}}(T):=\left(t_{11}\right)^{2 \sigma_{1}} \ldots\left(t_{r r}\right)^{2 \sigma_{r}}\left(T \in H_{\mathcal{V}}\right)$. Note that any one-dimensional representation $\chi$ of $H_{\mathcal{V}}$ is of the form $\chi_{\underline{\sigma}}$, so that $\chi$ is determined by the values on the subgroup $A_{\mathcal{V}} \subset H_{\mathcal{V}}$ consisting of diagonal matrices.

Let us give some examples. When $n_{1}=n_{2}=\cdots=n_{r}=1$ and $\mathcal{V}_{l k}=\mathbb{R}$ for all $1 \leq k<l \leq r$, the conditions (V1)-(V3) are clearly satisfied, and we have $\mathcal{Z}_{\mathcal{V}}=\operatorname{Sym}(r, \mathbb{R})$ and $\mathcal{P}_{\mathcal{V}}=\Pi_{r}$. For the case $r=3, n_{1}=n_{2}=n_{3}=1, \mathcal{V}_{21}=\{0\}$ and $\mathcal{V}_{31}=\mathcal{V}_{32}=\mathbb{R}$, the space $\mathcal{Z}_{\mathcal{V}}$ equals $\mathcal{Z}$ in (2.5).

Let us set $r=3, n_{1}=2, n_{2}=n_{3}=1$,

$$
\mathcal{V}_{21}=\left\{\left(\begin{array}{ll}
v & 0
\end{array}\right) ; v \in \mathbb{R}\right\}, \quad \mathcal{V}_{31}=\left\{\left(\begin{array}{ll}
0 & v
\end{array}\right) ; v \in \mathbb{R}\right\},
$$

and $\mathcal{V}_{32}=\{0\}$. Then we have

$$
\mathcal{Z}_{\mathcal{V}}=\left\{\left(\begin{array}{cccc}
y_{11} & 0 & y_{21} & 0  \tag{3.5}\\
0 & y_{11} & 0 & y_{31} \\
y_{21} & 0 & y_{22} & 0 \\
0 & y_{31} & 0 & y_{33}
\end{array}\right) ; y_{11}, y_{22}, y_{33}, y_{21}, y_{31} \in \mathbb{R}\right\}
$$

and

$$
\begin{align*}
\mathcal{P}_{\mathcal{V}} & =\left\{y \in \mathcal{Z}_{\mathcal{V}} ; y \text { is positive definite }\right\} \\
& =\left\{y \in \mathcal{Z}_{\mathcal{V}} ; y_{11}>0, y_{11} y_{22}-\left(y_{21}\right)^{2}>0, y_{11} y_{33}-\left(y_{31}\right)^{2}>0\right\} \tag{3.6}
\end{align*}
$$

which is exactly the Vinberg cone [24].

Set $r=2, n_{1}=m \geq 1, n_{2}=1$ and $\mathcal{V}_{21}=\operatorname{Mat}(1, m ; \mathbb{R})$. Then

$$
\begin{align*}
& \mathcal{Z}_{\mathcal{V}}=\left\{\left(\begin{array}{cccc}
y_{11} & & & v_{1} \\
& \ddots & & \vdots \\
& & y_{11} & v_{m} \\
v_{1} & \cdots & v_{m} & y_{22}
\end{array}\right) ; y_{11}, y_{22}, v_{1}, \ldots, v_{m} \in \mathbb{R}\right\} \\
& \mathcal{P}_{\mathcal{V}}=\left\{y \in \mathcal{Z}_{\mathcal{V}} ; y_{11}>0, y_{11} y_{22}-\left(v_{1}\right)^{2}-\cdots-\left(v_{m}\right)^{2}>0\right\} \tag{3.7}
\end{align*}
$$

so that we obtain the Lorentz cone of dimension $m+2$.
3.3. Basic quadratic maps. Let $W_{\mathcal{V}}^{i}(i=1, \ldots, r)$ be the subspace of $\operatorname{Mat}\left(N, n_{i} ; \mathbb{R}\right)$ consisting of matrices $x$ of the form

$$
x=\left(\begin{array}{c}
0_{n_{1}+\cdots+n_{i-1}, n_{i}} \\
X_{i i} \\
\vdots \\
X_{r i}
\end{array}\right) \quad\binom{X_{i i}=x_{i i} I_{n_{i}}, x_{i i} \in \mathbb{R}}{X_{l i} \in \mathcal{V}_{l i}(l=i+1, \ldots, r)}
$$

For example, when $\mathcal{Z}_{\mathcal{V}}$ is the one in (3.5), we have

$$
\begin{aligned}
& W_{\mathcal{V}}^{1}=\left\{\left(\begin{array}{cc}
x_{11} & 0 \\
0 & x_{11} \\
x_{21} & 0 \\
0 & x_{31}
\end{array}\right) ; x_{11}, x_{21}, x_{31} \in \mathbb{R}\right\}, \\
& W_{\mathcal{V}}^{2}=\left\{\left(\begin{array}{c}
0 \\
0 \\
x_{22} \\
0
\end{array}\right) ; x_{22} \in \mathbb{R}\right\}, \quad W_{\mathcal{V}}^{3}=\left\{\left(\begin{array}{c}
0 \\
0 \\
0 \\
x_{33}
\end{array}\right) ; x_{33} \in \mathbb{R}\right\},
\end{aligned}
$$

while for the case that $\mathcal{Z}_{\mathcal{V}}$ is the space $\mathcal{Z}$ in (2.5), we have

$$
\begin{aligned}
& W_{\mathcal{V}}^{1}=\left\{\left(\begin{array}{c}
x_{11} \\
0 \\
x_{31}
\end{array}\right) ; x_{11}, x_{31} \in \mathbb{R}\right\}, \quad W_{\mathcal{V}}^{2}=\left\{\left(\begin{array}{c}
0 \\
x_{22} \\
x_{32}
\end{array}\right) ; x_{22}, x_{32} \in \mathbb{R}\right\}, \\
& W_{\mathcal{V}}^{3}=\left\{\left(\begin{array}{c}
0 \\
0 \\
x_{33}
\end{array}\right) ; x_{33} \in \mathbb{R}\right\} .
\end{aligned}
$$

For $T \in H_{\mathcal{V}}$ and $x \in W_{\mathcal{V}}^{i}$, we see from (V1) that $T x \in W_{\mathcal{V}}^{i}$, which defines a representation $\tau_{i}: H_{\mathcal{V}} \rightarrow G L\left(W_{\mathcal{V}}^{i}\right)$. Since $W_{\mathcal{V}}^{i} \subset \operatorname{Mat}\left(N, n_{i} ; \mathbb{R}\right)$, we can consider the restriction $q_{\mathcal{V}}^{i}$ of the $\Pi_{N}$-positive quadratic map $q_{N, n_{i}}: \operatorname{Mat}\left(N, n_{i} ; \mathbb{R}\right) \rightarrow \operatorname{Sym}(N, \mathbb{R})$ in Example 2 to the space $W_{\mathcal{V}}^{i}$. Thanks to (V2) and (V3), we have $q_{\mathcal{V}}^{i}(x)=x^{\mathrm{t}} x \in \mathcal{Z}_{\mathcal{V}}$ for $x \in W_{\mathcal{V}}^{i}$. Then the quadratic map $q_{\mathcal{V}}^{i}: W_{\mathcal{V}}^{i} \ni x \mapsto x^{\mathrm{t}} x \in \mathcal{Z}_{\mathcal{V}}$ is $\mathcal{P}_{\mathcal{V}}$-positive. On
the other hand, we observe

$$
\begin{equation*}
q_{\mathcal{V}}^{i}\left(\tau_{i}(T) x\right)=(T x)^{\mathrm{t}}(T x)=\rho(T) q_{\mathcal{V}}^{i}(x) \quad\left(x \in W_{\mathcal{V}}^{i}, T \in H_{\mathcal{V}}\right) \tag{3.8}
\end{equation*}
$$

which yields the group homomorphism $H_{\mathcal{V}} \ni T \mapsto\left(\rho(T), \tau_{i}(T)\right) \in \operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q_{\mathcal{V}}^{i}\right)$. It follows that the quadratic map $q_{\mathcal{V}}^{i}$ is homogeneous. We call $q_{\mathcal{V}}^{1}, \ldots, q_{\mathcal{V}}^{r}$ basic quadratic maps for $\mathcal{P}_{\mathcal{V}}$. Recalling the inner product on $\mathcal{V}_{l k}$ given by (3.3), we define an inner product on the space $W_{\mathcal{V}}^{i}$ via the natural isomorphism

$$
\begin{equation*}
W_{\mathcal{V}}^{i} \simeq \mathbb{R} \oplus \sum_{l>i}^{\oplus} \mathcal{V}_{l i} \tag{3.9}
\end{equation*}
$$

Taking an orthonormal basis of $W_{\mathcal{V}}^{i}$ with respect to the inner product, we identify $W_{\mathcal{V}}^{i}$ with $\mathbb{R}^{m(i)}$, where $m(i):=\operatorname{dim} W_{\mathcal{V}}^{i}$. Then we consider the linear map $\phi_{\mathcal{V}}^{i}$ : $\mathcal{Z}_{\mathcal{V}} \equiv \mathcal{Z}_{\mathcal{V}}^{*} \rightarrow \operatorname{Sym}(m(i), \mathbb{R})$ associated to the quadratic map $q_{\mathcal{V}}^{i}$. Note that, if we write $n_{l i}$ for $\operatorname{dim} \mathcal{V}_{l i}$, we have $m(i)=1+\sum_{l>i} n_{l i}$.

Proposition 3.1. (i) One has

$$
\operatorname{det} \tau_{i}(T)=\chi_{\underline{m}(i) / 2}(T) \quad\left(T \in H_{\mathcal{V}}\right)
$$

where $\underline{m}(i)=\left(0, \cdots, 0,1, n_{i+1, i}, \ldots, n_{r i}\right) \in \mathbb{Z}^{r}$.
(ii) For $\eta=\rho(T)^{*} I_{N} \in \mathcal{P}_{\mathcal{V}}^{*}$ with $T \in H_{\mathcal{V}}$, one has

$$
\operatorname{det} \phi_{\mathcal{V}}^{i}(\eta)=\chi_{\underline{m}(i)}(T)
$$

In particular, $\left(t_{r r}\right)^{2}=\operatorname{det} \phi_{\mathcal{V}}^{r}(\eta)$.
(iii) For $1 \leq i<r$, there exist integers $c_{i+1, i}, \ldots, c_{r i}$ such that

$$
\begin{gathered}
\left(t_{i i}\right)^{2}=\operatorname{det} \phi_{\mathcal{V}}^{i}(\eta) \cdot\left(\operatorname{det} \phi_{\mathcal{V}}^{i+1}(\eta)\right)^{c_{i+1, i}} \cdots\left(\operatorname{det} \phi_{\mathcal{V}}^{r}(\eta)\right)^{c_{r i}} \\
\left(\eta=\rho^{*}(T) I_{N} \in \mathcal{P}_{\mathcal{V}}^{*}, T \in H_{\mathcal{V}}\right) .
\end{gathered}
$$

(iv) One has $\mathcal{P}_{\mathcal{V}}^{*}=\left\{\eta \in \mathcal{Z}_{\mathcal{V}} ; \operatorname{det} \phi_{\mathcal{V}}^{i}(\eta)>0(i=1, \ldots, r)\right\}$.

Proof. (i) Since $H_{\mathcal{V}} \ni T \mapsto \operatorname{det} \tau_{i}(T) \in \mathbb{C}^{\times}$is a one-dimensional representation, it is sufficient to check the equality for diagonal matrices $T \in A_{\mathcal{V}}$. In this case, the isomorphism (3.9) gives the eigenspace decomposition of $\tau_{i}(T)$, where $\mathbb{R}$ and $\mathcal{V}_{l i}$ correspond to the eigenvalues $t_{i i}$ and $t_{l l}$ respectively. Therefore we have $\operatorname{det} \tau_{i}(T)=t_{i i} \prod_{l>i} t_{l i}^{n_{l i}}=\chi_{\underline{m}(i) / 2}(T)$.
(ii) Thanks to (3.3), we have $(x \mid x)=\left\langle q_{\mathcal{V}}^{i}(x), I_{N}\right\rangle$, which implies $\phi_{\mathcal{V}}^{i}\left(I_{N}\right)=I_{m(i)}$. Thus we get $\operatorname{det} \phi_{\mathcal{V}}^{i}(\eta)=\operatorname{det} \phi_{\mathcal{V}}^{i}\left(\rho(T)^{*} I_{N}\right)=\left(\operatorname{det} \tau_{i}(T)\right)^{2}=\chi_{\underline{m}(i)}(T)$ by (2.22) and (i).
(iii) We see from (ii) that $\left(t_{i i}\right)^{2}=\operatorname{det} \phi_{\mathcal{V}}^{i}(\eta) \cdot\left(t_{i+1, i+1}\right)^{-2 n_{i+1, i}} \cdots\left(t_{r r}\right)^{-2 n_{r i}}$, whence we
can deduce (iii) recursively.
(iv) It is known ([23, Chapter 3, Section 3] and [6, Section 1]) that a homogeneous cone is described as the subset of the ambient vector space consisting of points at which all relatively invariant (appropriately normalized) functions are positive. Thus the assertion follows from (iii).

Example 4. Let $\mathcal{P}_{\mathcal{V}}$ be the Vinberg cone, that is, $\mathcal{Z}_{\mathcal{V}}$ is as in (3.5). Then we have

$$
\phi_{\mathcal{V}}^{1}(\eta)=\left(\begin{array}{ccc}
\eta_{11} & \eta_{21} & \eta_{31} \\
\eta_{21} & \eta_{22} & 0 \\
\eta_{31} & 0 & \eta_{33}
\end{array}\right), \quad \phi_{\mathcal{V}}^{2}(\eta)=\eta_{22}, \quad \phi_{\mathcal{V}}^{3}(\eta)=\eta_{33}
$$

for $\eta \in \mathcal{Z}_{\mathcal{V}}$. If $\eta=\rho^{*}(T) I_{4}$ for $T \in H_{\mathcal{V}}$, we have

$$
\left(t_{11}\right)^{2}\left(t_{22}\right)^{2}\left(t_{33}\right)^{2}=\operatorname{det} \phi_{\mathcal{V}}^{1}(\eta), \quad\left(t_{22}\right)^{2}=\eta_{22}, \quad\left(t_{33}\right)^{2}=\eta_{33}
$$

so that

$$
\left(t_{11}\right)^{2}=\frac{\operatorname{det} \phi_{\mathcal{V}}^{1}(\eta)}{\eta_{22} \eta_{33}}=\eta_{11}-\frac{\left(\eta_{21}\right)^{2}}{\eta_{22}}-\frac{\left(\eta_{31}\right)^{2}}{\eta_{33}} .
$$

On the other hand, we have by Proposition 3.1 (iv)

$$
\mathcal{P}_{\mathcal{V}}^{*}=\left\{\eta \in \mathcal{Z}_{\mathcal{V}} ; \eta_{11} \eta_{22} \eta_{33}-\eta_{33}\left(\eta_{21}\right)^{2}-\eta_{22}\left(\eta_{31}\right)^{2}>0, \eta_{22}>0, \eta_{33}>0\right\}
$$

Therefore, if $\mathcal{Z}$ is the space in (2.5), the linear isomorphism

$$
\iota: \mathcal{Z}_{\mathcal{V}} \ni\left(\begin{array}{cccc}
\eta_{11} & 0 & \eta_{21} & 0 \\
0 & \eta_{11} & 0 & \eta_{31} \\
\eta_{21} & 0 & \eta_{22} & 0 \\
0 & \eta_{31} & 0 & \eta_{33}
\end{array}\right) \mapsto\left(\begin{array}{ccc}
\eta_{33} & 0 & \eta_{31} \\
0 & \eta_{22} & \eta_{21} \\
\eta_{31} & \eta_{21} & \eta_{11}
\end{array}\right) \in \mathcal{Z}
$$

gives a bijection from $\mathcal{P}_{\mathcal{V}}^{*}$ onto $\mathcal{P}=\mathcal{Z} \cap \Pi_{3}$. The adjoint map $\iota^{*}: \mathcal{Z}^{*} \rightarrow \mathcal{Z}_{\mathcal{V}}^{*} \equiv \mathcal{Z}_{\mathcal{V}}$ gives the matrix realization of the Vinberg cone $\mathcal{Q}$ as the homogeneous cone $\mathcal{P}_{\mathcal{V}}$.
3.4. Standard quadratic maps and $H$-orbits in $\overline{\mathcal{P}_{\mathcal{V}}}$. We call the maps $q^{\underline{\mathcal{V}}}(\underline{\varepsilon} \in$ $\left.\{0,1\}^{r}, \underline{\varepsilon} \neq(0, \ldots, 0)\right)$ standard quadratic maps. They are of particular importance among the virtual quadratic maps $q^{\frac{s}{\mathcal{V}}}=\left(q_{\mathcal{V}}^{1}\right)^{\oplus s_{1}} \oplus \cdots \oplus\left(q_{\mathcal{V}}^{r}\right)^{\oplus s_{r}}\left(\underline{s}=\left(s_{1}, \ldots, s_{r}\right) \in\right.$ $\left.\mathbb{R}^{r}\right)$. Let us denote by $I(\underline{\varepsilon})$ the set $\left\{1 \leq i \leq r ; \underline{\varepsilon}_{i}=1\right\}$. We identify $q_{\mathcal{V}}^{\varepsilon}$ with a direct sum $\sum_{i \in I(\underline{\varepsilon})}^{\oplus} q_{\mathcal{V}}^{i}$ on the space $W_{\mathcal{V}}^{\varepsilon}:=\sum_{i \in I(\underline{\varepsilon})}^{\oplus} W_{\mathcal{V}}^{i}$. Recalling (3.9), we have the isomorphism $W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}} \simeq \sum_{i \in I(\underline{\varepsilon})}^{\oplus}\left(\mathbb{R} \oplus \sum_{l>i}^{\oplus} \mathcal{V}_{l i}\right)$, which enables us to describe $y=q^{\frac{\varepsilon}{\mathcal{V}}}(x) \in$ $\mathcal{Z}_{\mathcal{V}}\left(x \in W_{\mathcal{V}}^{\varepsilon}\right)$ as the matrix composed of the blocks

$$
\begin{equation*}
Y_{l k}=\sum_{i \leq k, i \in I(\underline{\underline{\varepsilon}})} X_{l i}{ }^{\mathrm{t}} X_{k i} \quad(1 \leq k \leq l \leq r), \tag{3.10}
\end{equation*}
$$

where $X_{i i}:=x_{i i} I_{n_{i}}$ for $i \in I(\underline{\varepsilon})$. For the case $l=k$, we have $Y_{k k}=y_{k k} I_{n_{k}}$ and

$$
\begin{equation*}
y_{k k}=\sum_{i \leq k, i \in I(\underline{\varepsilon})}\left\|X_{k i}\right\|^{2} \tag{3.11}
\end{equation*}
$$

thanks to (3.3), where we put $\left\|X_{i i}\right\|:=\left|x_{i i}\right|$. For each $x \in W_{\mathcal{V}}^{\underline{\varepsilon}}$, let $T_{x} \in \operatorname{Mat}(N, \mathbb{R})$ be a lower triangular matrix whose $(k, i)$-block component is $X_{k i}$ for $k>i$ with $i \in I(\varepsilon)$, and other components are zero. Then we have

$$
\begin{equation*}
q_{\mathcal{V}}^{\varepsilon}(x)=T_{x}{ }^{\mathrm{t}} T_{x} . \tag{3.12}
\end{equation*}
$$

For example, if $r=3$ and $\underline{\varepsilon}=(1,0,1)$, then an element $x$ of $W_{\mathcal{V}}^{\varepsilon}=W_{\mathcal{V}}^{1} \oplus W_{\mathcal{V}}^{3}$ is of the form

$$
x=\left(\begin{array}{l}
X_{11} \\
X_{21} \\
X_{31}
\end{array}\right) \oplus\left(\begin{array}{c}
0 \\
0 \\
X_{33}
\end{array}\right) \quad\binom{X_{11}=x_{11} I_{n_{1}}, X_{33}=x_{33} I_{n_{3}}}{x_{11}, x_{33} \in \mathbb{R}, X_{21} \in \mathcal{V}_{21}, X_{31} \in \mathcal{V}_{31}}
$$

and we have

$$
T_{x}=\left(\begin{array}{lll}
X_{11} & & \\
X_{21} & 0 & \\
X_{31} & 0 & X_{33}
\end{array}\right) \in \operatorname{Mat}(N, \mathbb{R})
$$

For $\underline{\varepsilon} \in\{0,1\}^{r}$, let $E_{\underline{\varepsilon}}$ be the element of $\mathcal{V}$ given by

$$
E_{\underline{\varepsilon}}:=\left(\begin{array}{ccc}
\varepsilon_{1} I_{n_{1}} & & \\
& \ddots & \\
& & \varepsilon_{r} I_{n_{r}}
\end{array}\right)
$$

and $\mathcal{O}_{\underline{\varepsilon}}$ the $H_{\mathcal{V}}$-orbit $\rho\left(H_{\mathcal{V}}\right) E_{\underline{\varepsilon}} \subset \mathcal{Z}_{\mathcal{V}}$ through $E_{\underline{\varepsilon}}$. In particular, the orbit $\mathcal{O}_{(0, \ldots, 0)}$ is the origin $\{0\}$, while $\mathcal{O}_{(1, \ldots, 1)}=\rho\left(H_{\mathcal{V}}\right) I_{N}$ equals the cone $\mathcal{P}_{\mathcal{V}}$. It is shown in [11, Theorem 3.5] that the $H_{\mathcal{V}}$-orbit decomposition of the closure $\overline{\mathcal{P}_{\mathcal{V}}}$ is given as

$$
\overline{\mathcal{P}_{\mathcal{V}}}=\bigsqcup_{\underline{\varepsilon} \in\{0,1\}^{r}} \mathcal{O}_{\underline{\varepsilon}}
$$

Proposition 3.2. If $\underline{\varepsilon} \neq(0, \ldots, 0)$, the image of the quadratic map $q \frac{\varepsilon}{\mathcal{V}}$ equals the closure $\overline{\mathcal{O}_{\underline{\varepsilon}}}$ of the orbit $\mathcal{O}_{\underline{\varepsilon}}$.

Proof. Let $W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}},+$ be the subset $\left\{x \in W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}} ; x_{i i}>0(i \in I(\underline{\varepsilon}))\right\}$ of $W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}}$. For $y=$ $\rho(T) E_{\underline{\varepsilon}}=T E_{\underline{\varepsilon}}{ }^{\mathrm{t}} T \in \mathcal{O}_{\underline{\varepsilon}}$ with $T \in H_{\mathcal{V}}$, we take a unique $x \in W_{\mathcal{V}}^{\boldsymbol{\varepsilon}},+$ for which $T_{x}$ equals $T E_{\underline{\varepsilon}}$. Then we have $y=q_{\mathcal{V}}^{\underline{\varepsilon}}(x)$ by (3.12). Conversely, for any $x \in W_{\mathcal{V}}^{\underline{\mathcal{\varepsilon}}},+$, we put $T:=T_{x}+I-E_{\underline{\varepsilon}} \in H_{\mathcal{V}}$ so that $q q_{\mathcal{V}}^{\varepsilon}(x)=T E_{\underline{\varepsilon}}{ }^{\mathrm{t}} T \in \mathcal{O}_{\underline{\varepsilon}}$. Therefore we obtain $\mathcal{O}_{\underline{\varepsilon}}=$ $q_{\mathcal{V}}^{\varepsilon}\left(W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}},+\right.$. On the other hand, putting $W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}}, *:=\left\{x \in W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}} ; x_{i i} \neq 0(i \in I(\underline{\varepsilon}))\right\}$, we see easily that $q_{\mathcal{V}}^{\underline{\varepsilon}}\left(W_{\mathcal{V}}^{\underline{\varepsilon}}, *\right)=q_{\mathcal{V}}^{\varepsilon}\left(W_{\mathcal{V}}^{\underline{\varepsilon}},+\right.$. Since $W_{\mathcal{V}}^{\underline{\mathcal{V}}, *}$ is an open dense subset of $W_{\mathcal{V}}^{\underline{\varepsilon}}$, the orbit $\mathcal{O}_{\underline{\varepsilon}}$ is dense in the image of the quadratic map $q_{\overline{\mathcal{V}}}^{\varepsilon}$, which is necessarily
closed. Indeed, introducing the projective imbedding $\iota_{V}$ of a vector space $V$ by $\iota_{V}: V \ni y \mapsto[1, y] \in \mathbf{P} V:=(\mathbb{R} \times V \backslash\{(0,0)\}) / \mathbb{R}^{\times}$, we can extend $q_{\mathcal{V}}^{\varepsilon}: W_{\mathcal{V}}^{\varepsilon} \rightarrow \mathcal{Z}_{\mathcal{V}}$ to the map $\tilde{q}_{\mathcal{\nu}}^{\varepsilon}: \mathbf{P} W_{\overline{\mathcal{V}}}^{\varepsilon} \ni[t, x] \mapsto\left[t^{2}, q_{\mathcal{\nu}}^{\varepsilon}(x)\right] \in \mathbf{P} \mathcal{Z}_{\mathcal{V}}$ because $q_{\overline{\mathcal{V}}}^{\varepsilon}(x) \neq 0$ for $x \neq 0$. The image $\tilde{q}_{\mathcal{V}}^{\underline{\varepsilon}}\left(\mathbf{P} W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}}\right)$ is compact, so that $q_{\mathcal{V}}^{\varepsilon}\left(W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}}\right)=\iota_{\mathcal{Z}}^{-1}\left(\tilde{q}_{\mathcal{V}}^{\varepsilon}\left(\mathbf{P} W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}}\right)\right)$ is closed.

Remark 3.3. In the proof of Proposition [3.2, we see that the quadratic map $q \frac{\varepsilon}{\mathcal{V}}$ gives a surjective map from $W_{\overline{\mathcal{V}}}^{\underline{\varepsilon},+}$ onto $\mathcal{O}_{\underline{\varepsilon}}$. The map is also one-to-one thanks to [11, Lemma 3.3 (ii)], while the map $q_{\overline{\mathcal{V}}}^{\varepsilon}: W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}, *} \rightarrow \mathcal{O}_{\underline{\varepsilon}}$ is $2^{\sharp I(\underline{\varepsilon})}$-to-one. Actually, a large part of the content of this section is presented in language of normal $j$-algebra in [11, Sections 3 and 4].

We define the representation $\tau_{\underline{\varepsilon}}: H_{\mathcal{V}} \rightarrow G L\left(W_{\mathcal{V}}^{\underline{\varepsilon}}\right)$ as the direct sum of the representations $\left(\tau_{i}, W_{\mathcal{V}}^{i}\right)$ of $H_{\mathcal{V}}$ for $i \in I(\underline{\varepsilon})$. Then we have by (3.8)

$$
\begin{equation*}
\left.q_{\mathcal{V}}^{\varepsilon}\left(\tau_{\underline{\varepsilon}}(T) x\right)=\rho(T) q\right) \quad\left(x \in W_{\mathcal{V}}^{\varepsilon}, T \in H_{\mathcal{V}}^{\varepsilon}\right) \tag{3.13}
\end{equation*}
$$

which implies that $q \frac{\varepsilon}{\mathcal{V}}$ is homogeneous.
The open set $W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}}+\subset W_{\overline{\mathcal{V}}}^{\varepsilon}$ is preserved by the action of $\tau_{\underline{\varepsilon}}\left(H_{\mathcal{V}}\right)$. We put

$$
R_{+}(\underline{\varepsilon}):=\left\{\underline{u}=\left(u_{1}, \ldots, u_{r}\right) \in \mathbb{R}^{r} ; u_{i}=0\left(\text { if } \varepsilon_{i}=0\right), \quad u_{i}>0\left(\text { if } \varepsilon_{i}=1\right)\right\} .
$$

For $\underline{u} \in R_{+}(\underline{\varepsilon})$, let $\mathcal{M}_{\underline{\underline{u}}}^{\underline{\varepsilon}}$ be the measure on $W_{\underline{\mathcal{V}}}^{\underline{\varepsilon}}+$ given by

$$
\begin{align*}
\mathcal{M}_{\underline{\underline{u}}}^{\varepsilon}(d x) & :=\prod_{i \in I(\underline{\underline{\varepsilon}})}\left\{\frac{2\left(x_{i i}\right)^{2 u_{i}-1} d x_{i i}}{\Gamma\left(u_{i}\right)} \cdot \prod_{l>i} \frac{d X_{l i}}{\pi^{n_{l i} / 2}}\right\} \\
& =\frac{\prod_{i \in I(\underline{\varepsilon})}\left(x_{i i}\right)^{2 u_{i}-1}}{\Gamma_{\underline{\varepsilon}}(\underline{u})} d x \quad\left(x \in W_{\mathcal{V}}^{\varepsilon,+}\right), \tag{3.14}
\end{align*}
$$

where $\Gamma_{\underline{\varepsilon}}(\underline{u})=\pi^{\operatorname{dim} W_{\bar{v}}^{\underline{\varepsilon}} / 2} \prod_{i \in I(\underline{\varepsilon})} \frac{\Gamma\left(u_{i}\right)}{2 \sqrt{\pi}}$. When $\underline{u}=\underline{\varepsilon} / 2$, the measure $\mathcal{M}_{\underline{\varepsilon} / 2}^{\underline{\varepsilon}}$ equals a constant multiple of the Lebesgue measure, that is, $\mathcal{M}_{\underline{\varepsilon} / 2}^{\underline{\varepsilon}}(d x)=2^{\sharp I(\underline{\varepsilon})} \pi^{-\operatorname{dim} W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}} / 2} d x$. We define $\underline{p}(\underline{\varepsilon}):=\left(p_{1}(\underline{\varepsilon}), p_{2}(\underline{\varepsilon}), \ldots, p_{r}(\underline{\varepsilon})\right)$ by

$$
p_{k}(\underline{\varepsilon}):=\sum_{l>k} \varepsilon_{l} n_{l k} .
$$

Lemma 3.4. (i) For a measurable set $A \subset W_{\overline{\mathcal{V}}}^{\underline{\varepsilon},+}$, one has

$$
\mathcal{M}_{\underline{\underline{u}}}^{\underline{\varepsilon}}\left(\tau_{\underline{\varepsilon}}(T) A\right)=\chi_{\underline{u}+\underline{p}(\underline{\varepsilon}) / 2}(T) \mathcal{M}_{\underline{\underline{u}}}^{\underline{\varepsilon}}(A) \quad\left(T \in H_{\mathcal{V}}\right) .
$$

(ii) One has

$$
\begin{equation*}
\int_{W_{\bar{v}}^{\varepsilon},+} e^{-\|x\|^{2}} \mathcal{M} \underline{\underline{u}}(d x)=1 \tag{3.15}
\end{equation*}
$$

Proof. (i) If $x^{\prime}=\tau_{\underline{\varepsilon}}(T) x \in W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}},+$ with $x \in W_{\overline{\mathcal{V}}}^{\underline{\underline{\varepsilon}},+}$, we have $x_{i i}^{\prime}=t_{i i} x_{i i}$ for $i \in I(\underline{\varepsilon})$. Thus

$$
\begin{equation*}
\prod_{i \in I(\underline{\varepsilon})}\left(x_{i i}^{\prime}\right)^{2 u_{i}-1}=\chi_{\underline{u}-\underline{\varepsilon} / 2}(T) \prod_{i \in I(\underline{\varepsilon})}\left(x_{i i}\right)^{2 u_{i}-1} \tag{3.16}
\end{equation*}
$$

On the other hand, we observe that

$$
d x^{\prime}=\left|\operatorname{det} \tau_{\underline{\varepsilon}}(T)\right| d x=\left(\prod_{i \in I(\underline{\varepsilon})} \operatorname{det} \tau_{i}(T)\right) d x
$$

and the last term equals $\left(\prod_{i \in I(\underline{\varepsilon})} \chi_{\underline{m}(i) / 2}(T)\right) d x$ by Proposition 3.1 (i). Since

$$
\sum_{i \in I(\varepsilon)} \underline{m}(i)=\sum_{i=1}^{r} \varepsilon_{i} \underline{m}(i)=\underline{\varepsilon}+\underline{p}(\underline{\varepsilon}),
$$

we have $d x^{\prime}=\chi_{\underline{\varepsilon} / 2+\underline{p}(\varepsilon) / 2}(T) d x$, which together with (3.16) implies (i).
(ii) By definition, we have

$$
\|x\|^{2}=\sum_{i \in I(\underline{\varepsilon})}\left\{\left(x_{i i}\right)^{2}+\sum_{l>i}\left\|X_{l i}\right\|^{2}\right\} .
$$

Thus the left-hand side of (3.15) equals

$$
\prod_{i \in I(\underline{\varepsilon})}\left\{\frac{2}{\Gamma\left(u_{i}\right)} \int_{0}^{\infty} e^{-\left(x_{i i}\right)^{2}}\left(x_{i i}\right)^{2 u_{i}-1} d x_{i i} \prod_{l>i} \int_{\mathcal{V}_{l i}} e^{-\left\|X_{l i}\right\|^{2}} \frac{d X_{l i}}{\pi^{n_{l i} / 2}}\right\}
$$

Therefore we obtain (3.15) from $\int_{\mathcal{V}_{l i}} e^{-\left\|X_{l i}\right\|^{2}} d X_{l i}=\pi^{n_{l i} / 2}$ and $\int_{0}^{\infty} e^{-\left(x_{i i}\right)^{2}}\left(x_{i i}\right)^{2 u_{i}-1} d x_{i i}=$ $\Gamma\left(u_{i}\right) / 2$.

Remark 3.5. Lemma 3.4 tells us that $e^{-\|x\|^{2}} \mathcal{M} \underline{\underline{u}} \bar{\varepsilon}(d x)$ is a probability measure on $W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}},+$. Actually, we see from the proof that if $X \underline{\underline{u}}$ is an $W_{\underline{\mathcal{V}}}^{\underline{\varepsilon},+}$-valued random variable with the law $e^{-\|x\|^{2}} \mathcal{M} \underline{\underline{u}}(d x)$, then its components are independent and satisfy $\sqrt{2} X_{l i}^{\underline{u}} \sim N\left(0, I_{n_{l i}}\right)$, and $\left(\sqrt{2} X_{i i}^{\underline{u}}\right)^{2} \sim \chi^{2}\left(2 u_{i}\right)$, where $\chi^{2}(u)$ denotes the chi-square law with the density $2^{-u} \Gamma(u / 2)^{-1} e^{-t / 2} t^{u-1}(t>0)$. We shall see later that any Wishart law associated to a homogeneous quadratic map is the image of this measure by an appropriate quadratic map.
3.5. Gindikin-Riesz distributions. For $\underline{\sigma}=\left(\sigma_{1}, \ldots, \sigma_{r}\right) \in \mathbb{C}^{r}$, we denote $\left(\sigma_{r}, \ldots, \sigma_{1}\right)$ by $\underline{\sigma}^{*}$. Let $\Delta_{\underline{\sigma}}^{*}$ be the function on the cone $\mathcal{P}_{\mathcal{V}}^{*}$ given by

$$
\begin{equation*}
\Delta_{\underline{\sigma}}^{*}\left(\rho^{*}(T) I_{N}\right)=\chi_{\underline{\sigma}^{*}}(T) \quad\left(T \in H_{\mathcal{V}}\right) \tag{3.17}
\end{equation*}
$$

By Proposition 3.1 (ii) and (iii), $\Delta_{\underline{\sigma}}^{*}(\eta)$ can be expressed as a product of powers of the polynomials $\operatorname{det} \phi_{\mathcal{V}}^{i}(\eta)$. Putting

$$
E_{\underline{\eta}}:=\left(\begin{array}{ccc}
\eta_{1} I_{n_{1}} & & \\
& \ddots & \\
& & \eta_{r} I_{n_{r}}
\end{array}\right) \in \mathcal{Z}_{\mathcal{V}}
$$

for $\underline{\eta}=\left(\eta_{1}, \ldots, \eta_{r}\right) \in \mathbb{R}_{>0}^{r}$, we have

$$
\begin{equation*}
\Delta_{\underline{\sigma}}^{*}\left(E_{\underline{\eta}}\right)=\left(\eta_{1}\right)^{\sigma_{r}}\left(\eta_{2}\right)^{\sigma_{r-1}} \ldots\left(\eta_{r}\right)^{\sigma_{1}} . \tag{3.18}
\end{equation*}
$$

For each $\underline{\sigma} \in \mathbb{C}^{r}$, Gindikin ([6], [7]) constructed a tempered distribution $\mathcal{R}_{\underline{\sigma}} \in \mathcal{S}^{\prime}\left(\mathcal{Z}_{\mathcal{V}}\right)$ whose Laplace transform $L_{\mathcal{R}_{\underline{\sigma}}}(\theta)=\mathcal{R}_{\underline{\sigma}}\left(e^{\langle y, \theta\rangle}\right)$ is given by

$$
\begin{equation*}
L_{\mathcal{R}_{\underline{\sigma}}}(\theta)=\Delta_{-\underline{\sigma}^{*}}^{*}(-\theta) \quad\left(\theta \in-\mathcal{P}_{\mathcal{V}}^{*}\right) \tag{3.19}
\end{equation*}
$$

We call $\mathcal{R}_{\underline{\sigma}}$ the Gindikin-Riesz distribution on the homogeneous cone $\mathcal{P}_{\mathcal{V}}$. The support of $\mathcal{R}_{\underline{\sigma}}$ is contained in $\overline{\mathcal{P}_{\mathcal{V}}}$, and $\mathcal{R}_{\underline{\sigma}}$ is relatively invariant under the action of $\rho\left(H_{\mathcal{V}}\right)$, that is,

$$
\begin{equation*}
\mathcal{R}_{\underline{\sigma}}(f \circ \rho(T))=\chi_{-\underline{\sigma}}(T) \mathcal{R}_{\underline{\sigma}}(f) \tag{3.20}
\end{equation*}
$$

for $T \in H_{\mathcal{V}}$ and $f \in \mathcal{S}\left(\mathcal{Z}_{\mathcal{V}}\right)$.
Proposition 3.6. For non-zero $\underline{\varepsilon} \in\{0,1\}^{r}$ and $\underline{u} \in R_{+}(\underline{\varepsilon})$, put $\underline{\sigma}:=\underline{u}+\underline{p}(\underline{\varepsilon}) / 2$. Then $\mathcal{R}_{\underline{\boldsymbol{\sigma}}}$ is the image of $\mathcal{M}_{\underline{\underline{u}}}^{\underline{\varepsilon}}$ by the standard quadratic map $q \frac{\bar{\nu}}{\underline{\varepsilon}}$.

Proof. By (3.19), it is sufficient to show that

$$
\int_{W_{\underline{v}}^{\varepsilon},+} e^{\left\langle\frac{q_{\bar{\nu}}^{\varepsilon}}{\nu}(x), \theta\right\rangle} \mathcal{M}_{\underline{u}}^{\varepsilon}(d x)=\Delta_{-\underline{\sigma}^{*}}^{*}(-\theta)
$$

for $\theta \in-\mathcal{P}_{\mathcal{V}}^{*}$. Take $T \in H_{\mathcal{V}}$ for which $\theta=-\rho^{*}(T) I_{N}$. Then the left-hand side is

$$
\int_{W_{\bar{v}}^{\varepsilon},+} e^{-\left\langle q_{\hat{\nu}}^{\varepsilon}(x), \rho^{*}(T) I_{N}\right\rangle} \mathcal{M}_{\underline{u}}^{\varepsilon}(d x)=\int_{W_{\bar{v}}^{\varepsilon},+} e^{-\left\langle q_{\underline{\mathcal{V}}}^{\varepsilon}\left(\tau_{\underline{\varepsilon}}(T) x\right), I_{N}\right\rangle} \mathcal{M}_{\underline{u}}^{\frac{\varepsilon}{\underline{\varepsilon}}}(d x)
$$

by (3.13), and it is equal to

$$
\chi_{-(\underline{u}+\underline{p}(\underline{\varepsilon}) / 2)}(T) \int_{W_{\underline{\underline{v}}}^{\varepsilon},+} e^{-\left\langle q_{\hat{\nu}}^{\varepsilon}(x), I_{N}\right\rangle} \mathcal{M}_{\underline{\underline{u}}}^{\underline{\varepsilon}}(d x)
$$

by Lemma 3.4 (i). Since $\left\langle q_{\mathcal{\nu}}^{\varepsilon}(x), I_{N}\right\rangle=\|x\|^{2}$ by (3.11), we see from Lemma 3.4 (ii) that

$$
\left.\int_{W_{\bar{v}}^{\varepsilon},+} e^{\left\langle q_{\underline{\mathcal{V}}}^{\varepsilon}(x), \theta\right\rangle} \mathcal{M}_{\underline{\underline{u}}}^{\varepsilon}(d x)=\chi_{-(\underline{u}+\underline{p}}(\underline{\varepsilon}) / 2\right)(T)=\Delta_{-\underline{\sigma}^{*}}^{*}(-\theta) .
$$

We set

$$
\begin{aligned}
\Xi(\underline{\varepsilon}) & :=\left\{\underline{\sigma}=\underline{u}+\underline{p}(\underline{\varepsilon}) / 2 ; \underline{u} \in R_{+}(\underline{\varepsilon})\right\} \\
& =\left\{\underline{\sigma} \in \mathbb{R}^{r} ; \sigma_{i}=p_{i}(\underline{\varepsilon}) / 2\left(\text { if } \varepsilon_{i}=0\right), \quad \sigma_{i}>p_{i}(\underline{\varepsilon}) / 2\left(\text { if } \varepsilon_{i}=1\right)\right\} .
\end{aligned}
$$

If $\underline{\varepsilon} \neq(0, \ldots, 0)$ and $\underline{\sigma} \in \Xi(\underline{\varepsilon})$, then $\mathcal{R}_{\underline{\sigma}}$ is a positive measure on the orbit $\mathcal{O}_{\underline{\varepsilon}}$ by Proposition 3.6. For the case $\varepsilon=(0, \ldots, 0)$, we have $\Xi(0, \ldots, 0)=\{(0, \ldots, 0)\}$ and $\mathcal{R}_{(0, \ldots, 0)}$ is the Dirac measure at the origin $\{0\}$. It is proven in 11] that they exhaust all the cases that $\mathcal{R}_{\underline{\sigma}}$ is a positive measure.

Theorem 3.7 ([11, Theorem 6.2]). The Gindikin-Riesz distribution $\mathcal{R}_{\underline{\sigma}}$ is a positive measure if and only if $\underline{\sigma} \in \Xi:=\bigsqcup_{\underline{\varepsilon} \in\{0,1\}^{r}} \Xi(\varepsilon)$. Moreover, if $\underline{\sigma} \in \Xi(\underline{\varepsilon})$, then $\mathcal{R}_{\underline{\sigma}}$ is a measure on $\mathcal{O}_{\underline{\varepsilon}}$.

The parameter set $\Xi$ is also described as

$$
\Xi=\left\{\sum_{i=1}^{r} \varepsilon_{i}\left(0, \cdots, 0, u_{i}, n_{i+1, i} / 2, \ldots, n_{r i} / 2\right) ; \varepsilon_{i} \in\{0,1\}, u_{i}>0(i=1, \ldots, r)\right\} .
$$

3.6. Riesz measures and Gindikin-Riesz distributions. Let us investigate a relation of the Riesz measures $\mu_{q}$ associated to homogeneous $\mathcal{P}_{\mathcal{V}}$-positive quadratic maps $q$ and the Gindikin-Riesz distributions on $\mathcal{P}_{\mathcal{V}}$.

Proposition 3.8. For $i=1, \ldots$, r, one has $\mu_{q_{\dot{\nu}}^{i}}=\pi^{m(i) / 2} \mathcal{R}_{\underline{m}(i)}$.
Proof. ¿From Lemma 2.1 and Proposition 3.1 (ii), we have

$$
\begin{equation*}
L_{\mu_{q_{\mathcal{V}}^{i}}}(\theta)=\pi^{m(i) / 2} \operatorname{det} \phi_{\mathcal{V}}^{i}(-\theta)^{-1 / 2}=\pi^{m(i) / 2} \Delta_{-\underline{m}(i)^{*} / 2}^{*}(-\theta) \quad\left(\theta \in-\mathcal{P}_{\mathcal{V}}^{*}\right) \tag{3.21}
\end{equation*}
$$

which implies the statement.

Assume that there exists the Riesz measure $\mu_{q_{V}^{s}}$ associated to a virtual quadratic $\operatorname{map} q_{\mathcal{V}}^{s}=\left(q_{\mathcal{V}}^{1}\right)^{\oplus s_{1}} \oplus \cdots \oplus\left(q_{\mathcal{V}}^{r}\right)^{\oplus s_{r}}$. By (2.12) and (3.21), we have for $\theta \in-\mathcal{P}_{\mathcal{V}}^{*}$

$$
L_{\mu_{q \underline{s}}}(\theta)=\prod_{i=1}^{r}\left(\pi^{m(i) / 2} \Delta_{-\underline{m}(i)^{*} / 2}^{*}(-\theta)\right)^{s_{i}}
$$

We put

$$
\begin{equation*}
\underline{\sigma}:=\sum_{i=1}^{r} s_{i} \underline{m}(i) / 2=\frac{1}{2} \sum_{i=1}^{r} s_{i}\left(0, \ldots, 0,1, n_{i+1, i}, \cdots, n_{r i}\right) . \tag{3.22}
\end{equation*}
$$

Then the equality above can be rewritten as

$$
L_{\mu_{q_{\mathcal{V}}^{s}}^{\underline{s}}}(\theta)=\pi^{|\underline{\underline{q}}|} \Delta_{-\underline{\sigma}^{*}}^{*}(-\theta),
$$

where $|\underline{\sigma}|:=\sigma_{1}+\cdots+\sigma_{r}$. Thus $\mu_{q_{\bar{\nu}}}$ equals $\pi^{|\underline{\sigma}|} \mathcal{R}_{\underline{\sigma}}$, so that $\underline{\sigma}$ belongs to $\Xi(\underline{\varepsilon})$ for some $\underline{\varepsilon} \in\{0,1\}^{r}$ owing to Theorem 3.7. The converse argument is also valid. Therefore we obtain

Theorem 3.9. For a virtual quadratic map $q_{\mathcal{V}}^{s}=\left(q_{\mathcal{V}}^{1}\right)^{\oplus s_{1}} \oplus \cdots \oplus\left(q_{\mathcal{V}}^{r}\right)^{\oplus s_{r}}$, there exists the associated Riesz measure $\mu_{q_{\mathcal{V}}}$ if and only if $\underline{\sigma}:=\sum_{i=1}^{r} s_{i} \underline{m}(i) / 2$ belongs to $\Xi$. In this case $\mu_{q^{\underline{\nu}}}=\pi^{|\underline{\sigma}|} \mathcal{R}_{\underline{\sigma}}$, and there exist $\underline{\varepsilon} \in\{0,1\}^{r}$ and $\underline{u} \in R_{+}(\underline{\varepsilon})$ for which $\underline{\sigma}=\underline{u}+\underline{p}(\underline{\varepsilon}) / 2$. If $\underline{\varepsilon} \neq(0, \ldots, 0), \mu_{q \underline{\mathcal{v}}}$ is the image of the measure $\pi^{|\sigma|} \mathcal{M}_{\underline{\underline{u}}}^{\underline{\varepsilon}}$ on $W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}}$ by the standard quadratic map $q^{\frac{\varepsilon}{\nu}}$.

Let $q: \mathbb{R}^{m} \rightarrow \mathcal{Z}_{\mathcal{V}}$ be any homogeneous $\mathcal{P}_{\mathcal{V}}$-positive quadratic map. As is noted in Section 3.1, the group $p r_{1}(\operatorname{Aut}(\mathcal{P} \mathcal{V}, q))$ acts on the cone $\mathcal{P}_{\mathcal{V}}$ transitively. Assume first that $\operatorname{pr}_{1}\left(\operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q\right)\right)$ contains $\rho\left(H_{\mathcal{V}}\right)$. Then the polynomial $\operatorname{det} \phi_{q}(\eta)$ is relatively invariant under the action of $\rho\left(H_{\mathcal{V}}\right)$ by (2.23). Namely, for each $T \in H_{\mathcal{V}}$, there exists $c_{T}>0$ such that $\operatorname{det} \phi_{q}\left(\rho^{*}(T) \eta\right)=c_{T} \operatorname{det} \phi_{q}(\eta)\left(\eta \in\left(\mathbb{R}^{n}\right)^{*}\right)$. It is easy to see that the correspondence $H_{\mathcal{V}} \ni T \mapsto c_{T} \in \mathbb{R}_{>0}$ is a one-dimensional representation, so that we have $c_{T}=\chi_{\underline{m}}(T)$ for some $\underline{m} \in \mathbb{R}^{r}$. Thus $\operatorname{det} \phi_{q}\left(\rho^{*}(T) I_{n}\right)=\chi_{\underline{\underline{m}}}(T) \operatorname{det} \phi_{q}\left(I_{N}\right)$ for $T \in H_{\mathcal{V}}$, which means that $\operatorname{det} \phi_{q}(\eta)=C \Delta_{\underline{m}^{*}}^{*}(\eta)$ for $\eta \in \mathcal{P}_{\mathcal{V}}^{*}$, where $C:=\operatorname{det} \phi_{q}\left(I_{N}\right)$. By (3.18), we have

$$
\operatorname{det} \phi_{q}\left(E_{\underline{\eta}}\right)=C\left(\eta_{1}\right)^{m_{1}} \ldots\left(\eta_{r}\right)^{m_{r}}
$$

which gives a practical way to determine $m_{i}$. Indeed, we see from this formula that $m_{i}$ are non-negative integers. Comparing the degrees of both sides, we obtain $m=m_{1}+\cdots+m_{r}$. Similarly to Proposition 3.8, we have

$$
\begin{equation*}
\mu_{q}=C^{-1 / 2} \pi^{m / 2} \mathcal{R}_{\underline{m} / 2} \tag{3.23}
\end{equation*}
$$

Let us consider the virtual quadratic map $q^{\oplus s}$. The associated Riesz measure exists if and only if $s \underline{m} / 2 \in \Xi$, and in this case

$$
\mu_{q^{\oplus s}}=C^{-s / 2} \pi^{s m / 2} \mathcal{R}_{s \underline{m} / 2}
$$

As for the general case, we have the following result.
Proposition 3.10. Let $q: \mathbb{R}^{m} \rightarrow \mathcal{Z}_{\mathcal{V}}$ be a homogeneous $\mathcal{P}_{\mathcal{V}}$-positive quadratic map. Then there exist $g_{0} \in G\left(\mathcal{P}_{\mathcal{V}}\right), \underline{m} \in \mathbb{Z}^{r}$, and $C>0$ for which

$$
\operatorname{det} \phi_{q}\left(\left(g_{0}^{-1}\right)^{*} \eta\right)=C \Delta_{\underline{m}^{*}}^{*}(\eta) \quad\left(\eta \in \mathcal{P}_{\mathcal{V}}\right)
$$

The Riesz measure $\mu_{q^{\oplus s}}$ associated to the virtual quadratic map $q^{\oplus s}$ exists if and only if $s \underline{m} / 2 \in \Xi$. In this case, $\mu_{q^{\oplus s}}$ equals the image of $C^{-s / 2} \pi^{s m / 2} \mathcal{R}_{s \underline{m} / 2}$ by $g_{0}$.

Proof. We note that $\operatorname{pr}_{1}\left(\operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q\right)\right)$ acts on the cone $\mathcal{P}_{\mathcal{V}}$ transitively, and that the identity component of $p r_{1}\left(\operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q\right)\right)$ equals the identity component of an algebraic group (cf. [15, Theorem 2]). It follows that an Iwasawa subgroup (maximal connected split solvable subgroup) $\mathcal{H}$ of $\operatorname{pr}_{1}\left(\operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q\right)\right)$ acts on $\mathcal{P}_{\mathcal{V}}$ simply transitively $([23$, Chapter 1$])$. Since $\mathcal{H}$ is also an Iwasawa subgroup of $G\left(\mathcal{P}_{\mathcal{V}}\right)$, it is conjugate to another Iwasawa subgroup $\rho\left(H_{\mathcal{V}}\right) \subset G\left(\mathcal{P}_{\mathcal{V}}\right)$. Namely, there exists $g_{0} \in G\left(\mathcal{P}_{\mathcal{V}}\right)$ for which $g_{0}^{-1} \mathcal{H} g_{0}=\rho\left(H_{\mathcal{V}}\right)$. Let $q^{\prime}$ be the $\mathcal{P}$-positive quadratic $\operatorname{map} g_{0}^{-1} \circ q: \mathbb{R}^{m} \rightarrow \mathcal{Z}_{\mathcal{V}}$. We have $\phi_{q^{\prime}}(\eta)=\phi_{q}\left(\left(g_{0}^{-1}\right)^{*} \eta\right)$ for $\eta \in\left(\mathbb{R}^{n}\right)^{*}$ because ${ }^{\mathrm{t}} x \phi_{q^{\prime}}(\eta) x=\left\langle q^{\prime}(x), \eta\right\rangle=\left\langle q(x),\left(g_{0}^{-1}\right)^{*} \eta\right\rangle={ }^{\mathrm{t}} x \phi_{q}\left(\left(g_{0}^{-1}\right)^{*} \eta\right) x$ for $x \in \mathbb{R}^{m}$. It is easy to see that

$$
\operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q^{\prime}\right)=\left\{\left(g_{0}^{-1} g_{1} g_{0}, g_{2}\right) \in G L\left(\mathcal{Z}_{\mathcal{V}}\right) \times G L\left(\mathbb{R}^{m}\right) ;\left(g_{1}, g_{2}\right) \in \operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q\right)\right\}
$$

Then $\operatorname{pr}_{1}\left(\operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q^{\prime}\right)\right)=g_{0}^{-1} \operatorname{pr}_{1}\left(\operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q\right)\right) g_{0} \supset g_{0}^{-1} \mathcal{H} g_{0}=\rho\left(H_{\mathcal{V}}\right)$. Thus we can apply the argument preceding Proposition 3.10 for $q^{\prime}$, so that we have

$$
\operatorname{det} \phi_{q}\left(\left(g_{0}^{-1}\right)^{*} \eta\right)=\operatorname{det} \phi_{q^{\prime}}(\eta)=C \Delta_{\underline{m}^{*}}^{*}(\eta)
$$

with some $C>0$ and $\underline{m} \in \mathbb{Z}^{r}$. Moreover $\mu_{\left(q^{\prime}\right) \oplus s}$ equals $C^{-s / 2} \pi^{s m / 2} \mathcal{R}_{s \underline{m} / 2}$ if $s \underline{m} / 2 \in \Xi$. Since $q^{\oplus s}=g_{0} \circ\left(q^{\prime}\right)^{\oplus s}$, we get the last statement from Proposition 2.12.

Proposition 3.10 states that the Riesz measure $\mu_{q}$ associated to a homogeneous $q$ is equal to some Gindikin-Riesz distribution up to a linear transform on $G\left(\mathcal{P}_{\mathcal{V}}\right)$. On the other hand, Theorem 3.9 tells us that if a Gindikin-Riesz distribution is a positive measure, then it equals a Riesz measure associated to the virtual sum $q^{\frac{s}{\mathcal{V}}}=\left(q_{\mathcal{V}}^{1}\right)^{\oplus s_{1}} \oplus \cdots \oplus\left(q_{\mathcal{V}}^{r}\right)^{\oplus s_{r}}$ of basic quadratic maps up to a constant multiple. For example, let us recall the homogeneous $\Pi_{r}$-positive quadratic map $q^{I}: R^{I} \rightarrow \operatorname{Sym}(r, \mathbb{R})$ with $I \subset\{1, \ldots, r\}$ and the permutation matrix $w_{0} \in \mathfrak{S}_{r} \subset G L(r, \mathbb{R})$ in Section 3.1. Putting $g_{0}:=\rho\left(w_{0}\right)$, we have $q^{I}=g_{0} \circ q^{\{r-k+1, \ldots, r\}}(k:=\sharp I)$, while $q^{\{r-k+1, \ldots, r\}}$ is exactly the basic quadratic map $q_{\mathcal{V}}^{r-k+1}$ for $\mathcal{Z}_{\mathcal{V}}=\operatorname{Sym}(r, \mathbb{R})$. Therefore, the Riesz measure $\mu_{\left(q^{I}\right)^{\oplus s}}$ exists if and only if $(0, \ldots, 0, \underbrace{s / 2, \ldots, s / 2}_{k}) \in \Xi$, that is, $s \in\{0,1, \ldots, k-1\} \cup(k-1,+\infty)$. In this case, $\mu_{\left(q^{I}\right)^{\oplus}}$ equals the image of $\pi^{s / 2} \mathcal{R}_{(0, \ldots, 0, s / 2, \ldots, s / 2)}=\mu_{\left(q_{\nu}^{r-k+1}\right)^{\oplus s}}$ by $g_{0}$.

Let $q_{1}: \mathbb{R}^{m_{1}} \rightarrow \mathcal{Z}_{\mathcal{V}}$ and $q_{2}: \mathbb{R}^{m_{2}} \rightarrow \mathcal{Z}_{\mathcal{V}}$ be two homogeneous $\mathcal{P}_{\mathcal{V}}$-positive quadratic maps. As we have seen in Section 3.1, the direct sum $q_{1} \oplus q_{2}$ is not necessarily homogeneous. Let us assume that the group $p r_{1}\left(\operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q_{1}\right)\right) \cap p r_{1}\left(\operatorname{Aut}\left(\mathcal{P}, q_{2}\right)\right)$ acts
on $\mathcal{P}_{\mathcal{V}}$ transitively. In this case, we see easily that $q_{1} \oplus q_{2}$ is homogeneous. As in Proposition 3.10, we can take $g_{0} \in G\left(\mathcal{P}_{\mathcal{V}}\right)$ for which $g_{0} \rho\left(H_{\mathcal{V}}\right) g_{0}^{-1} \subset p r_{1}\left(\operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q_{1}\right)\right) \cap$ $\operatorname{pr}_{1}\left(\operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q_{2}\right)\right)$. Then we have $\operatorname{det} \phi_{q_{1}}\left(\left(g_{0}^{-1}\right)^{*} \eta\right)=C_{1} \Delta_{\underline{m}^{\prime}}^{*}(\eta)$ and $\operatorname{det} \phi_{q_{2}}\left(\left(g_{0}^{-1}\right)^{*} \eta\right)=$ $C_{2} \Delta_{\underline{m}^{\prime \prime}}^{*}(\eta)$ for $\eta \in\left(\mathbb{R}^{n}\right)^{*}$ with some $C_{1}, C_{2}>0$ and $\underline{m}^{\prime}, \underline{m}^{\prime \prime} \in \mathbb{Z}^{r}$. Now we consider a virtual quadratic map $q=q_{1}^{\oplus s_{1}} \oplus q_{2}^{\oplus s_{2}}$. We see that the associated Riesz measure $\mu_{q}$ exists if and only if $s_{1} \underline{m}^{\prime}+s_{2} \underline{m}^{\prime \prime} / 2 \in \Xi$, and in this case, $\mu_{q}$ is the image of $C_{1}^{-s_{1} / 2} C_{2}^{-s_{2} / 2} \pi^{\left(s_{1} m_{1}+s_{2} m_{2}\right) / 2} \mathcal{R}_{\left(s_{1} \underline{m}^{\prime}+s_{2} \underline{m}^{\prime \prime}\right) / 2}$ by $g_{0}$. Obviously, the same argument is valid for general quadratic maps $q=q_{1}^{\oplus s_{1}} \oplus q_{2}^{\oplus s_{2}} \oplus \cdots \oplus q_{t}^{\oplus s_{t}}$.
3.7. Bartlett decomposition of the Wishart laws. Let $q: \mathbb{R}^{m} \rightarrow \mathcal{Z}_{\mathcal{V}}$ be a homogeneous $\mathcal{P}_{\mathcal{V}}$-positive quadratic map. Then the Wishart law $\gamma_{q, \theta}\left(\theta \in-\mathcal{P}_{\mathcal{V}}^{*}\right)$ is the image of the normal law $N\left(0, \phi(-\theta)^{-1}\right)$ on the vector space $\mathbb{R}^{m}$ by the quadratic map $q / 2$, see Remark 2.3.

However, this description of the Wishart law does not permit us to determine its support in general. In this section, we shall give another construction of the Wishart random matrices, which is a generalization of the Bartlett decomposition ([2], [20, Theorem 3.2.14]) and has the advantage of controlling the support of the underlying Wishart law. Moreover, the result is valid for virtual quadratic maps.

First we consider the virtual quadratic map $q_{\mathcal{V}}^{s}=\left(q_{\mathcal{V}}^{1}\right)^{\oplus s_{1}} \oplus \cdots \oplus\left(q_{\mathcal{V}}^{r}\right)^{\oplus s_{r}}$ whose associated Riesz measure $\mu_{q \bar{\nu}}$ exists. Then $\underline{\sigma}=\sum_{i=1}^{r} s_{i} \underline{m}(i) / 2$ belongs to $\Xi$ and we have $\mu_{q \frac{s}{\mathcal{V}}}=\pi^{|\underline{\sigma}|} \mathcal{R}_{\underline{\sigma}}$ by Theorem [3.9. Moreover, we have $L_{\mu_{q \frac{s}{\nu}}}(\theta)=\pi^{|\underline{\sigma}|} \Delta_{-\underline{\sigma}^{*}}^{*}(-\theta)$. Therefore we obtain from (2.8) that

$$
\begin{equation*}
\gamma_{q_{\bar{v}}^{s}, \theta}(d y)=e^{\langle y, \theta\rangle} \Delta_{\sigma^{*}}^{*}(-\theta) \mathcal{R}_{\underline{\sigma}}(d y) \quad\left(y \in \mathbb{R}^{n}\right) \tag{3.24}
\end{equation*}
$$

We remark that distributions of this type are considered in [10] for the case when $\mathcal{P}_{\mathcal{V}}$ is a symmetric cone. Assume that $\gamma_{q_{\mathcal{V}}, \theta}$ is not the Dirac measure. Then $\underline{\sigma} \neq(0, \ldots, 0)$, so that we can take a non-zero $\underline{\varepsilon} \in\{0,1\}^{r}$ and $\underline{u} \in R_{+}(\underline{\varepsilon})$ for which $\underline{\sigma}=\underline{u}+\underline{p}(\underline{\varepsilon}) / 2$. Recall the standard quadratic map $q^{\varepsilon}: W_{\overline{\mathcal{V}}}^{\underline{\varepsilon}} \rightarrow \mathcal{Z}_{\mathcal{V}}$ and the subset $W_{\mathcal{V}}^{\underline{\varepsilon}},+\subset W_{\mathcal{V}}^{\varepsilon}$ introduced in Section 3.4. As noted in (3.12), each element $x \in W_{\mathcal{V}}^{\underline{\varepsilon}},+$ is identified with a lower triangular matrix $T_{x}$ for which $q_{\bar{\nu}}^{\varepsilon}(x)=T_{x}{ }^{\mathrm{t}} T_{x}$. Thus, the $W_{\mathcal{V}}^{\underline{\varepsilon}},+$-valued random variable $X^{\underline{u}}$ in the following theorem can be regarded as a triangular random matrix, similarly to the Bartlett decomposition of the classical Wishart distribution.

Theorem 3.11. Let $\underline{\sigma}=\underline{u}+\underline{p}(\underline{\varepsilon}) / 2$ and $X^{\underline{u}}$ be an $W_{\overline{\mathcal{V}}}^{\underline{\mathcal{E}},+}$-valued random variable whose components are independent and satisfy $\left(X_{i i}^{\underline{u}}\right)^{2} \sim \chi^{2}\left(2 u_{i}\right)$ and $X_{l i}^{\underline{u}} \sim N\left(0, I_{n_{l i}}\right)$
for $i \in I(\underline{\varepsilon})$ and $l>i$.
(i) The Wishart law $\gamma_{q \overline{\mathcal{V}}},-I_{N}$ is the law of $Y=q_{\mathcal{V}}^{\varepsilon}\left(X^{\underline{u}}\right) / 2$ and is supported by $\mathcal{O}_{\underline{\varepsilon}}$.
(ii) For $\theta=-\rho(T)^{*} I_{N} \in-\mathcal{P}_{\mathcal{V}}^{*}$ with $T \in H_{\mathcal{V}}$, the Wishart law $\gamma_{q_{\mathcal{V}}, \theta}$ is the law of $Y^{\prime}=\rho(T)^{-1} \circ q_{\mathcal{\nu}}^{\underline{\varepsilon}}\left(X^{\underline{u}}\right) / 2$ and is supported by $\mathcal{O}_{\underline{\varepsilon}}$.

Proof. For a measurable function $f$ on $\mathcal{Z}_{\mathcal{V}}$, we see from (3.24) and Proposition 3.6 that
$\int_{\mathcal{Z}_{\mathcal{V}}} f(y) \gamma_{q^{\frac{s}{\mathcal{V}}},-I_{N}}(d y)=\int_{\mathcal{Z}_{\mathcal{V}}} f(y) e^{-\left\langle y, I_{N}\right\rangle} \mathcal{R}_{\underline{\boldsymbol{\sigma}}}(d y)=\int_{W_{\underline{\mathcal{V}}}^{\underline{\varepsilon}},+} f\left(q^{\frac{\mathcal{V}}{\varepsilon}}(x)\right) e^{-\left\langle q_{\mathcal{V}}^{\varepsilon}(x), I_{N}\right\rangle} \mathcal{M}_{\underline{\underline{u}}}^{\underline{\varepsilon}}(d x)$.
Since $\left\langle q^{\frac{\varepsilon}{\mathcal{L}}}(x), I_{N}\right\rangle=\|x\|^{2}$, by the change of variable of $x$ by $x / \sqrt{2}$, we rewrite the last term as

$$
\int_{W_{\underline{v}}^{\varepsilon},+} f\left(q^{\frac{\mathcal{V}}{\mathcal{V}}}(x) / 2\right) e^{-\|x\|^{2} / 2} \mathcal{M}_{\underline{\underline{u}}}^{\varepsilon}(d x / 2) .
$$

Keeping the Remark 3.5 in mind, we see that the law of the random variable $X^{\underline{u}}$ is $e^{-\|x\|^{2} / 2} \mathcal{M} \underline{\underline{u}} \boldsymbol{\varepsilon}(d x / 2)$. Hence (i) holds. To show (ii), it suffices to check that $\gamma_{q \bar{q}, ~}$, is the image of $\gamma_{q \bar{\nu}}^{s},-I_{N}$ by $\rho(T)^{-1}$. Since $\gamma_{q \overline{\mathcal{V}},-I_{N}}(d y)=e^{-\left\langle y, I_{N}\right\rangle} \mathcal{R}_{\underline{\sigma}}(d y)$, we have

$$
\gamma_{q \mathcal{V},-I_{N}}(\rho(T) d y)=e^{-\left\langle\rho(T) y, I_{N}\right\rangle} \mathcal{R}_{\underline{\boldsymbol{\sigma}}}(\rho(T) d y)=e^{\langle y, \theta\rangle} \chi_{\underline{\boldsymbol{\sigma}}}(T) \mathcal{R}_{\underline{\boldsymbol{\sigma}}}(d y)
$$

by (3.20). Therefore (3.24) together with (3.17) leads us to the assertion (ii).

Now we consider the Wishart distribution $\gamma_{q^{\oplus s}, \theta}$, where $q$ is a general homogeneous $\mathcal{P}_{\mathcal{V}}$-positive quadratic map. First we show a refinement of the first part of Proposition 3.10.

Lemma 3.12. Let $q: \mathbb{R}^{m} \rightarrow \mathcal{Z}_{\mathcal{V}}$ be a homogeneous $\mathcal{P}_{\mathcal{V}}$-positive quadratic map, and $\theta$ an element of $-\mathcal{P}_{\mathcal{V}}^{*}$. Then there exist $g_{0} \in G\left(\mathcal{P}_{\mathcal{V}}\right), \underline{m} \in \mathbb{Z}^{r}$ and $C>0$ for which

$$
\begin{equation*}
\operatorname{det} \phi_{q}\left(\left(g_{0}^{-1}\right)^{*} \eta\right)=C \Delta_{\underline{m}}^{*}(\eta) \quad\left(\eta \in \mathcal{P}_{\mathcal{V}}^{*}\right) \tag{3.25}
\end{equation*}
$$

and $g_{0}^{*} \theta=-I_{N}$.
Proof. It is shown in the proof of Proposition 3.10 that there exists $a_{0} \in G(\mathcal{P} \mathcal{V})$ for which $a_{0} \rho\left(H_{\mathcal{V}}\right) a_{0}^{-1} \subset \operatorname{pr}_{1}\left(\operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q\right)\right)$. Since $-a_{0}^{*} \theta \in \mathcal{P}_{\mathcal{V}}^{*}$, we take $T_{0} \in H_{\mathcal{V}}$ for which $\rho^{*}\left(T_{0}\right) I_{N}=-a_{0}^{*} \theta$. Put $g_{0}:=a_{0} \rho\left(T_{0}\right)^{-1} \in G\left(\mathcal{P}_{\mathcal{V}}\right)$. Then $g_{0} \rho\left(H_{\mathcal{V}}\right) g_{0}^{-1} \subset$ $\operatorname{pr}_{1}\left(\operatorname{Aut}\left(\mathcal{P}_{\mathcal{V}}, q\right)\right)$ and $g_{0}^{*} \theta=\rho^{*}\left(T_{0}\right)^{-1} a_{0}^{*} \theta=-I_{N}$. Similarly to Proposition 3.10, we see that $g_{0}$ together with an appropriate $\underline{m}$ and $C>0$ satisfies the required properties.

Assume that there exists the Riesz measure $\mu_{q^{\oplus s}}$ associated to a virtual quadratic map $q^{\oplus s}$, and that $\mu_{q^{\oplus s}}$ is not the Dirac measure. Then we can take non-zero
$\underline{\varepsilon} \in\{0,1\}^{r}$ and $\underline{u} \in R_{+}(\underline{\varepsilon})$ such that $s \underline{m} / 2=\underline{u}+\underline{p}(\underline{\varepsilon}) / 2$ as in Theorem 3.9. Using these data together with $g_{0}$ in Lemma 3.12, we obtain the Bartlett decomposition of the Wishart distribution $\gamma_{q^{\oplus s}, \theta}$.

Theorem 3.13. Let $s \underline{m} / 2=\underline{u}+\underline{p}(\underline{\varepsilon}) / 2$ and $X^{\underline{u}}$ be the $W_{\underline{\mathcal{V}}}^{\underline{\varepsilon}}$, -valued random variable in Theorem 3.11. Then the Wishart law $\gamma_{q^{\oplus s}, \theta}$ is the law of $Y=g_{0} \circ q_{\mathcal{V}}^{\varepsilon}\left(X^{\underline{u}}\right) / 2$.

Proof. Put $q^{\prime}:=g_{0}^{-1} \circ q^{\oplus s}$. As is seen in the proof of Proposition 3.10, the Riesz measure $\mu_{\left(q^{\prime}\right)^{\oplus s}}$ equals $C^{-s / 2} \pi^{s m / 2} \mathcal{R}_{s \underline{m} / 2}$. Thus, similarly to the proof of Theorem 3.11 (i), we see that $\gamma_{\left(q^{\prime}\right) \oplus s,-I_{N}}(d y)=e^{-\left\langle y, I_{N}\right\rangle} \mathcal{R}_{s \underline{m} / 2}(d y)$, and that $\gamma_{\left(q^{\prime}\right) \oplus s,-I_{N}}$ is the law of $Y=q_{\mathcal{V}}^{\varepsilon}\left(X^{\underline{u}}\right) / 2$. Since $q^{\oplus s}=g_{0} \circ\left(q^{\prime}\right)^{\oplus s}$, Theorem 3.13 follows from Proposition 2.12.

We have seen that Riesz measures and Wishart laws associated to a homogeneous quadratic map are obtained (up to linear transforms as in Proposition 3.10 and Theorem 3.13) as the ones associated to a virtual quadratic map $q^{\frac{s}{\mathcal{V}}}=\left(q_{\mathcal{V}}^{1}\right)^{\oplus s_{1}} \oplus$ $\cdots \oplus\left(q_{\mathcal{\nu}}^{r}\right)^{\oplus s_{r}}$, that is, a virtual sum of basic quadratic maps. However, it does not mean that every homogeneous quadratic map is equal to a direct sum of basic quadratic maps. The structure of homogeneous quadratic maps is more rich than the maps generated by basic quadratic maps. Let us study the following example.

Example 5. Let $\operatorname{Herm}(2, \mathbb{C})$ be the vector space of Hermitian matrices of size 2, and $\Omega \subset \operatorname{Herm}(2, \mathbb{C})$ the subset of positive definite matrices. Then we see that

$$
\Omega=\left\{\left(\begin{array}{cc}
y_{1} & y_{3}+i y_{4} \\
y_{3}-i y_{4} & y_{2}
\end{array}\right) ; y_{1}>0, y_{1} y_{2}-\left(y_{3}\right)^{2}-\left(y_{4}\right)^{2}>0\right\}
$$

so that $\Omega$ is the 4 -dimensional Lorentz cone. Recalling (3.7), we have the linear isomorphism

$$
\iota: \operatorname{Herm}(2, \mathbb{C}) \ni\left(\begin{array}{cc}
y_{1} & y_{3}-i y_{4} \\
y_{3}+i y_{4} & y_{2}
\end{array}\right) \mapsto\left(\begin{array}{ccc}
y_{1} & 0 & y_{3} \\
0 & y_{1} & y_{4} \\
y_{3} & y_{4} & y_{2}
\end{array}\right) \in \mathcal{Z}_{\mathcal{V}}
$$

which gives a matrix realization of $\Omega$. Let us consider the quadratic map $\tilde{q}: \mathbb{C}^{2} \ni$ $z \mapsto z^{\mathrm{t}} \bar{z} \in \operatorname{Herm}(2, \mathbb{C})$, which is clearly $\Omega$-positive. We have a group homomorphism

$$
\operatorname{GL}(2, \mathbb{C}) \ni A \mapsto(\tilde{\rho}(A), A) \in \operatorname{Aut}(\Omega, \tilde{q}),
$$

where $\tilde{\rho}(A) \in G L(\operatorname{Herm}(2, \mathbb{C}))$ is defined by $\tilde{\rho}(A)(Z):=A Z^{\mathrm{t}} \bar{A}(Z \in \operatorname{Herm}(2, \mathbb{C}))$. Since $\tilde{\rho}(G L(2, \mathbb{C}))$ acts on $\Omega$ transitively, the quadratic map $\tilde{q}$ is homogeneous. Keeping the natural isomorphism $\mathbb{C}^{2} \simeq \mathbb{R}^{4}$ in mind, we define the quadratic map
$q: \mathbb{R}^{4} \rightarrow \mathcal{Z}_{\mathcal{V}}$ by

$$
\begin{aligned}
q(x) & :=\iota \circ \tilde{q}\binom{x_{1}+i x_{2}}{x_{3}+i x_{4}} \\
& =\iota\left(\begin{array}{ccc}
\left(x_{1}\right)^{2}+\left(x_{2}\right)^{2} & \left(x_{1} x_{3}+x_{2} x_{4}\right)-i\left(x_{1} x_{4}-x_{2} x_{3}\right) \\
\left(x_{1} x_{3}+x_{2} x_{4}\right)+i\left(x_{1} x_{4}-x_{2} x_{3}\right) & \left(x_{3}\right)^{2}+\left(x_{4}\right)^{2}
\end{array}\right) \\
& =\left(\begin{array}{ccc}
\left(x_{1}\right)^{2}+\left(x_{2}\right)^{2} & 0 & x_{1} x_{3}+x_{2} x_{4} \\
0 & \left(x_{1}\right)^{2}+\left(x_{2}\right)^{2} & x_{1} x_{4}-x_{2} x_{3} \\
x_{1} x_{3}+x_{2} x_{4} & x_{1} x_{4}-x_{2} x_{3} & \left(x_{3}\right)^{2}+\left(x_{4}\right)^{2}
\end{array}\right) .
\end{aligned}
$$

Then we have

$$
\phi_{q}(\eta)=\left(\begin{array}{cccc}
\eta_{1} & 0 & \eta_{3} & \eta_{4} \\
0 & \eta_{1} & -\eta_{4} & \eta_{3} \\
\eta_{3} & -\eta_{4} & \eta_{2} & 0 \\
\eta_{4} & \eta_{3} & 0 & \eta_{2}
\end{array}\right)
$$

for $\eta \in \mathcal{Z}_{\mathcal{V}}$. It is easily checked that the map $\phi_{q} \circ \iota: \operatorname{Herm}(2, \mathbb{C}) \rightarrow \operatorname{Sym}(4, \mathbb{R})$ is a Jordan algebra representation. For $\eta \in \Omega^{*}$ we have

$$
\begin{equation*}
L_{\mu_{q}}(-\eta)=\pi^{2}\left(\operatorname{det} \phi_{q}(\eta)\right)^{-1 / 2}=\pi^{2}\left(\eta_{1} \eta_{2}-\left(\eta_{3}\right)^{2}-\left(\eta_{4}\right)^{2}\right)^{-1} \tag{3.26}
\end{equation*}
$$

by Lemma 2.1. On the other hand, the basic quadratic maps $q_{\mathcal{V}}^{i}: W_{\mathcal{V}}^{i} \rightarrow \mathcal{Z}_{\mathcal{V}}(i=$ 1,2 ) are given by

$$
q_{\mathcal{V}}^{1}\left(\begin{array}{cc}
x_{1} & 0 \\
0 & x_{1} \\
x_{3} & x_{4}
\end{array}\right)=\left(\begin{array}{ccc}
\left(x_{1}\right)^{2} & 0 & x_{1} x_{3} \\
0 & \left(x_{1}\right)^{2} & x_{1} x_{4} \\
x_{1} x_{3} & x_{1} x_{4} & \left(x_{3}\right)^{2}+\left(x_{4}\right)^{2}
\end{array}\right), \quad q_{\mathcal{V}}^{2}\left(\begin{array}{c}
0 \\
0 \\
x_{2}
\end{array}\right)=\left(\begin{array}{ccc}
0 & 0 & 0 \\
0 & 0 & 0 \\
0 & 0 & \left(x_{2}\right)^{2}
\end{array}\right)
$$

so that we have for $\eta \in \mathcal{Z}_{\mathcal{V}}^{*}$

$$
\phi_{\mathcal{V}}^{1}(\eta)=\left(\begin{array}{ccc}
\eta_{1} & \eta_{3} & \eta_{4} \\
\eta_{3} & \eta_{2} & 0 \\
\eta_{4} & 0 & \eta_{2}
\end{array}\right), \quad \phi_{\mathcal{V}}^{2}(\eta)=\eta_{2}
$$

Thus we obtain

$$
\begin{equation*}
L_{\mu_{q_{\mathcal{V}}^{1}}}(-\eta)=\pi^{3 / 2}\left(\eta_{2}\right)^{-1 / 2}\left(\eta_{1} \eta_{2}-\left(\eta_{3}\right)^{2}-\left(\eta_{4}\right)^{2}\right)^{-1 / 2}, \quad L_{\mu_{q_{\mathcal{V}}^{2}}}(-\eta)=\pi^{1 / 2}\left(\eta_{2}\right)^{-1 / 2} \tag{3.27}
\end{equation*}
$$

for $\eta \in \mathcal{P}_{\mathcal{V}}^{*}$. Comparing (3.26) and (3.27), we see that $\mu_{q}=\mu_{\left(q_{\mathcal{V}}^{1} \oplus^{\oplus 2} \oplus\left(q_{\mathcal{V}}^{2}\right)^{\oplus(-2)}\right.}$, whereas the quadratic map $q$ is by no means equal to the virtual quadratic map $\left(q_{\mathcal{V}}^{1}\right)^{\oplus 2} \oplus$ $\left(q_{\mathcal{V}}^{2}\right)^{\oplus(-2)}$. We see also from (3.26) and (3.27) that $\mu_{q \oplus\left(q_{\mathcal{V}}^{2}\right)^{2}}=\mu_{\left(q_{\mathcal{V}}^{1}\right)^{\oplus} \text {, whereas the }}$ two (true) quadratic maps $q \oplus\left(q_{\mathcal{V}}^{2}\right)^{\oplus 2}$ and $\left(q_{\mathcal{V}}^{1}\right)^{\oplus 2}$ do not coincide even up to linear transforms $g_{0} \in G(\Omega)$, as the domains of these maps are different.

Therefore two different quadratic maps may correspond to the same Riesz measure.
3.8. Density function for the non-singular case. Since the orbit $\mathcal{O}_{\underline{\varepsilon}}=$ $\rho(H) E_{\underline{\varepsilon}}$ is contained in the boundary $\partial \mathcal{P}_{\mathcal{V}}$ of the homogeneous cone $\mathcal{P}_{\mathcal{V}}$ unless $\varepsilon=$ $(1, \ldots, 1)$, the Gindikin-Riesz distribution $\mathcal{R}_{\underline{\sigma}}$ is a singular measure for $\underline{\sigma} \in \Xi(\underline{\varepsilon})$ with $\underline{\varepsilon} \neq(1, \ldots, 1)$ thanks to Proposition 3.6. On the other hand, if $\underline{\sigma} \in \Xi(1, \ldots, 1)$, that is,

$$
\sigma_{i}>p_{i} / 2 \quad(i=1, \ldots, r),
$$

where $p_{i}:=p_{i}(1, \ldots, 1)=\sum_{l>i} n_{l i}$, then the Gindikin-Riesz distribution is an absolutely continuous measure with respect to the Lebesgue measure, and the density function is given explicitly in [6] as follows.

Noting that the group $H_{\mathcal{V}}$ acts on $\mathcal{P}_{\mathcal{V}}$ simply transitively, we define the function $\Delta_{\underline{\sigma}}: \mathcal{P} \rightarrow \mathbb{C}^{\times}$for $\underline{\sigma}=\left(\sigma_{1}, \ldots, \sigma_{r}\right) \in \mathbb{C}^{r}$ by $\Delta_{\underline{\sigma}}\left(\rho(T) I_{N}\right):=\chi_{\underline{\sigma}}(T)\left(T \in H_{\mathcal{V}}\right)$. For $y=\rho(T) I_{N}=T^{\mathrm{t}} T \in \mathcal{P}_{\mathcal{V}}$, we can express $\Delta_{\underline{\boldsymbol{q}}}(y)$ as a product of powers of principal minors of $y$ (cf. [4, p. 122]). Define $\underline{d}=\left(d_{1}, \ldots, d_{r}\right) \in \mathbb{Z}^{r} / 2$ by $d_{k}:=$ $1+\left(\sum_{l>k} n_{l k}+\sum_{i<k} n_{k i}\right) / 2$. Then $\Delta_{-\underline{d}}(y) d y$ gives a $G\left(\mathcal{P}_{\mathcal{V}}\right)$-invariant measure on $\mathcal{P}$ ([6, Proposition 2.2]). Take $\underline{\sigma} \in \Xi(1, \ldots, 1)$. We see from [6, Theorem 2.1] that the integral

$$
\Gamma_{\mathcal{P}_{\mathcal{V}}}(\underline{\sigma}):=\int_{\mathcal{P}_{\mathcal{V}}} e^{-\left\langle y, I_{N}\right\rangle} \Delta_{\underline{\sigma}-\underline{d}}(y) d y
$$

converges and equals $\pi^{\left(\operatorname{dim} \mathcal{Z}_{\mathcal{V}}-r\right) / 2} \prod_{i=1}^{r} \Gamma\left(\sigma_{i}-p_{i} / 2\right)$. By [6, Proposition 2.3], we see that

$$
\begin{equation*}
\mathcal{R}_{\underline{\sigma}}(d y)=\frac{\Delta_{\underline{\sigma}-\underline{d}}(y)}{\Gamma_{\mathcal{P}}(\underline{s})} d y \quad\left(y \in \mathcal{P}_{\mathcal{V}}\right) . \tag{3.28}
\end{equation*}
$$

Owing to (3.24) and (3.28), we conclude the following proposition.
Proposition 3.14. Let $q^{\frac{s}{\nu}}=\left(q_{\nu}^{1}\right)^{\oplus s_{1}} \oplus \cdots \oplus\left(q_{\mathcal{V}}^{r}\right)^{\oplus s_{r}}$ be the virtual quadratic map such that $\underline{\sigma}=\sum_{i=1}^{r} s_{i} \underline{\underline{m}}(i) / 2$ belongs to $\Xi(1, \ldots, 1)$, that is, $\sigma_{i}>p_{i} / 2(i=1, \ldots, r)$. Then one has

$$
\begin{equation*}
\gamma_{q_{\overline{\mathcal{V}}}^{s}, \theta}(d y)=\frac{e^{\langle y, \theta\rangle} \Delta_{\sigma^{*}}^{*}(-\theta) \Delta_{\underline{\sigma}-\underline{d}}(y)}{\Gamma_{\mathcal{P}_{\mathcal{V}}}(\underline{s})} d y \quad\left(y \in \mathcal{P}_{\mathcal{V}}\right) \tag{3.29}
\end{equation*}
$$

Note that the formula (3.29) served as a definition of a Wishart law in [1].
Example 6. Let $\mathcal{Z}_{\mathcal{V}}$ be the space defined in (3.5). If $y=\rho(T) I_{N}=T^{\mathrm{t}} T \in \mathcal{P}_{\mathcal{V}}$ with $T \in H_{\mathcal{V}}$, then we see easily that

$$
y_{11}=\left(t_{11}\right)^{2}, \quad y_{11} y_{22}-\left(y_{21}\right)^{2}=\left(t_{11}\right)^{2}\left(t_{22}\right)^{2}, \quad y_{11} y_{33}-\left(y_{31}\right)^{2}=\left(t_{11}\right)^{2}\left(t_{33}\right)^{2}
$$

so that $\Delta_{\underline{\sigma}}(y)=\left(t_{11}\right)^{2 \sigma_{1}}\left(t_{22}\right)^{2 \sigma_{2}}\left(t_{33}\right)^{2 \sigma_{3}}$ equals

$$
\left(y_{11}\right)^{\sigma_{1}-\sigma_{2}-\sigma_{3}}\left(y_{11} y_{22}-\left(y_{21}\right)^{2}\right)^{\sigma_{2}}\left(y_{11} y_{33}-\left(y_{31}\right)^{2}\right)^{\sigma_{3}} .
$$

On the other hand, we have $\left(p_{1}, p_{2}, p_{3}\right)=(0,1,1)$ and $\left(d_{1}, d_{2}, d_{3}\right)=(2,3 / 2,3 / 2)$. Thus we have by (3.28)

$$
\mathcal{R}_{\left(\sigma_{1}, \sigma_{2}, \sigma_{3}\right)}(d y)=\frac{\left(y_{11}\right)^{\sigma_{1}-\sigma_{2}-\sigma_{3}+1}\left(y_{11} y_{22}-\left(y_{21}\right)^{2}\right)^{\sigma_{2}-3 / 2}\left(y_{11} y_{33}-\left(y_{31}\right)^{2}\right)^{\sigma_{3}-3 / 2}}{\pi \Gamma\left(\sigma_{1}\right) \Gamma\left(\sigma_{2}-1 / 2\right) \Gamma\left(\sigma_{3}-1 / 2\right)} d y
$$

if $\sigma_{1}>0, \sigma_{2}>1 / 2$ and $\sigma_{3}>1 / 2$. Let us consider the Wishart laws associated to the virtual quadratic $\operatorname{map}\left(q_{\mathcal{V}}^{1}\right)^{\oplus s}$, where $q_{\mathcal{V}}^{1}: W_{\mathcal{V}}^{1} \rightarrow \mathcal{Z}_{\mathcal{V}}$ is the basic quadratic map. Since $\underline{m}(1)=(1,1,1)$, we observe that $s \underline{m}(1) / 2 \in \Xi$ if and only if $s \in$ $\{0,1\} \cup(1,+\infty)$. If $s=0$, the associated Wishart law is the Dirac measure. If $s=1$, the associated Wishart law $\gamma_{q_{\mathcal{V}}^{1}, \theta}\left(\theta \in-\mathcal{P}_{\mathcal{V}}^{*}\right)$ is described as the image of the normal law $N\left(0, \phi_{\mathcal{V}}^{1}(-\theta)^{-1}\right)$ on $W_{\mathcal{V}}^{1} \equiv \mathbb{R}^{3}$ by the quadratic map $q_{\mathcal{V}}^{1} / 2$, where $\phi_{\mathcal{V}}^{1}(-\theta)$ is given in Example 4 after Proposition 3.1.

If $s>1$, then $s \underline{m}(1) / 2=(s / 2, s / 2, s / 2)$ belongs to $\Xi(1,1,1)$. Since

$$
\Delta_{s \underline{m}(1)^{*} / 2}^{*}(\eta)=\operatorname{det} \phi_{\mathcal{V}}^{1}(\eta)^{s / 2}=\left(\eta_{11} \eta_{22} \eta_{33}-\eta_{33}\left(\eta_{21}\right)^{2}-\eta_{22}\left(\eta_{31}\right)^{2}\right)^{s / 2} \quad\left(\eta \in \mathcal{P}_{\mathcal{V}}^{*}\right)
$$

we have for $\theta=-\eta \in-\mathcal{P}_{\mathcal{V}}^{*}$

$$
\begin{aligned}
\gamma_{\left(q_{\mathcal{V}}^{1}\right)^{s}, \theta}(d y)= & \frac{e^{-\langle y, \eta\rangle}\left(\eta_{11} \eta_{22} \eta_{33}-\eta_{33}\left(\eta_{21}\right)^{2}-\eta_{22}\left(\eta_{31}\right)^{2}\right)^{s / 2}}{\pi \Gamma(s / 2) \Gamma((s-1) / 2) \Gamma((s-1) / 2)} \times \\
& \times\left(y_{11}\right)^{1-s / 2}\left(y_{11} y_{22}-\left(y_{21}\right)^{2}\right)^{(s-3) / 2}\left(y_{11} y_{33}-\left(y_{31}\right)^{2}\right)^{(s-3) / 2} d y \quad\left(y \in \mathcal{P}_{\mathcal{V}}\right)
\end{aligned}
$$

by Proposition 3.14

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