# THE CLASS OF DISTRIBUTIONS OF PERIODIC ORNSTEIN-UHLENBECK PROCESSES <br> DRIVEN BY LÉVY PROCESSES 

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#### Abstract

The class $I(c)$ of stationary distributions of periodic Ornstein-Uhlenbeck processes with parameter $c$ driven by Lévy processes is analyzed. A characterization of $I(c)$ analogous to a well-known characterization of the selfdecomposable distributions is given. The relations between $I(c)$ for varying values of $c$ and the relations with the class of selfdecomposable distributions and with the nested classes $L_{m}$ are discussed.


## 1. INTRODUCTION

Let $\left\{Z_{t}\right\}_{t \in[0,1]}$ be a Lévy process restricted to $t \in[0,1]$ with values in $\mathbb{R}^{d}$. Let $c \in \mathbb{R} \backslash\{0\}$. Consider the Langevin equation

$$
\begin{equation*}
\mathrm{d} X_{t}=-c X_{t} \mathrm{~d} t+\mathrm{d} Z_{t}, \quad t \in[0,1] \tag{1.1}
\end{equation*}
$$

with boundary condition

$$
\begin{equation*}
X_{0}=X_{1} \quad \text { a.s. } \tag{1.2}
\end{equation*}
$$

This has a unique solution

$$
\begin{equation*}
X_{t}=e^{-c t} X_{0}+e^{-c t} \int_{0}^{t} e^{c s} \mathrm{~d} Z_{s}, \quad t \in[0,1] \tag{1.3}
\end{equation*}
$$

with

$$
\begin{equation*}
X_{0}=X_{1}=\frac{1}{e^{c}-1} \int_{0}^{1} e^{c s} \mathrm{~d} Z_{s} \tag{1.4}
\end{equation*}
$$

All equalities involving stochastic integrals are understood in the sense "almost surely". We call $\left\{X_{t}\right\}_{t \in[0,1]}$ the periodic Ornstein-Uhlenbeck process with parameter $c$ and background driving Lévy process $\left\{Z_{t}\right\}_{t \in[0,1]}$. The pathwise periodic extension $\left\{\widetilde{X}_{t}\right\}_{t \in \mathbb{R}}$ of $\left\{X_{t}\right\}_{t \in[0,1]}$ by $\widetilde{X}_{t}=X_{t}$ for $t \in[0,1]$ and $\widetilde{X}_{t+n}=\widetilde{X}_{t}$ for $t \in \mathbb{R}$ and integers $n$ is

[^0]a stationary process, that is, $\left\{\widetilde{X}_{t}\right\}_{t \in \mathbb{R}}$ and $\left\{\widetilde{X}_{t+s}\right\}_{t \in \mathbb{R}}$ have an identical system of finite-dimensional distributions for any $s \in \mathbb{R}$. In particular, $\mathcal{L}\left(X_{t}\right)=\mathcal{L}\left(X_{1}\right)=$ $\mathcal{L}\left(\frac{1}{e^{c}-1} \int_{0}^{1} e^{c s} \mathrm{~d} Z_{s}\right)$ for any $t \in[0,1]$, where we denote the distribution of a random vector $Y$ by $\mathcal{L}(Y)$. We call this distribution the stationary distribution of $\left\{X_{t}\right\}_{t \in[0,1]}$. These are processes introduced and studied by Pedersen [8]. The facts above were shown in [8] in the case where $d=1$ and $c>0$, and proofs in the general case can be given in the same way (see also [2] and [10] for the equivalence of (1.1) and (1.3)). The Gaussian periodic Ornstein-Uhlenbeck processes appear in Kwakernaak [6], and Grenander [3] and others use these processes in relation to statistical shape analysis; see [8].

In this paper we are interested in the analysis of the class $I(c)$ of stationary distributions of periodic Ornstein-Uhlenbeck processes with parameter $c$ where the background driving Lévy processes $\left\{Z_{t}\right\}_{t \in[0,1]}$ are ranging over all Lévy processes on $\mathbb{R}^{d}$. This continues the study initiated by Pedersen [8]. A usual Ornstein-Uhlenbeck process with parameter $c>0$ and background driving Lévy process $\left\{Z_{t}\right\}_{t \in[0, \infty)}$ is a solution of equation (1.1) with $t \in[0,1]$ replaced by $t \in[0, \infty)$ under the condition that $X_{0}$ is independent of the process $\left\{Z_{t}\right\}_{t \in[0, \infty)}$. It is determined uniquely by $X_{0}$ and expressed by (1.3) with $t \in[0,1]$ replaced by $t \in[0, \infty)$. It satisfies $\mathcal{L}\left(X_{t}\right)=\mathcal{L}\left(X_{0}\right)$ for $t \in[0, \infty)$ if and only if $\mathcal{L}\left(Z_{1}\right)$ has finite $\log$-moment and $\mathcal{L}\left(X_{0}\right)=\mathcal{L}\left(\int_{0}^{\infty} e^{-c s} \mathrm{~d} Z_{s}\right)$. In this case the stationary distribution $\mu=\mathcal{L}\left(X_{0}\right)$ is selfdecomposable. Conversely, any selfdecomposable distribution $\mu$ is expressed in this way with a Lévy process $\left\{Z_{t}\right\}_{t \in[0, \infty)}$ having finite log-moment. See [4] or [12]. The correspondence between $\mu$ and $\mathcal{L}\left(Z_{1}\right)$ is one-to-one for fixed $c>0$. Thus the class of stationary distributions of usual Ornstein-Uhlenbeck processes does not depend on the parameter $c$ and coincides with the class $S D\left(\mathbb{R}^{d}\right)$ of selfdecomposable distributions on $\mathbb{R}^{d}$. However, we will see that the class $I(c)$ delicately depends on the parameter $c$.

In the next section we will characterize $I(c)$ in terms of properties of the Lévy measures of the distributions in this class. The mapping $\Phi_{c}$ of $\mu_{0}=\mathcal{L}\left(Z_{1}\right)$ to $\mu=\Phi_{c}\left(\mu_{0}\right)=\mathcal{L}\left(\int_{0}^{1} e^{c s} \mathrm{~d} Z_{s}\right)$ is shown to be one-to-one and a homeomorphism. In Section 3 we will study the relations between $I(c)$ for varying values of $c$ and the relations between $I(c)$ and $S D\left(\mathbb{R}^{d}\right)$. It will be shown that $I(c) \supsetneqq S D\left(\mathbb{R}^{d}\right)$ and that $\bigcap_{n=1}^{\infty} I\left(c_{n}\right)=S D\left(\mathbb{R}^{d}\right)$ whenever $c_{n} \uparrow \infty$. One way of looking at the dependence of $I(c)$ on $c$ is to study, given $\mu$ in the class $I D\left(\mathbb{R}^{d}\right)$ of infinitely divisible distributions on $\mathbb{R}^{d}$, the set $Q(\mu)=\{c \geqslant 0: \mu \in I(c)\}$, where we define $I(0)=I D\left(\mathbb{R}^{d}\right)$. We will show
that these sets $Q(\mu)$ have a rich variety. In Section 4 we consider the relations to other classes of distributions. The decreasing sequence of classes $L_{m}, m=0,1, \ldots, \infty$, beginning with $L_{0}=S D\left(\mathbb{R}^{d}\right)$ is considered. These classes are known to be intimately connected with the mapping $\widetilde{\Phi}_{-c}$ of $\mu_{0}=\mathcal{L}\left(Z_{1}\right)$ to $\mu=\widetilde{\Phi}_{-c}\left(\mu_{0}\right)=\mathcal{L}\left(\int_{0}^{\infty} e^{-c t} \mathrm{~d} Z_{t}\right)$, but the relations with $\Phi_{c}$ are more subtle; we will show that $L_{m} \varsubsetneqq \Phi_{c}\left(L_{m-1}\right) \varsubsetneqq L_{m-1}$ for any finite $m$.

Stochastic integrals $\int_{0}^{t} f(s) \mathrm{d} Z_{s}$ in this paper are defined as limits in probability from the case of step functions $f$, as in [7], [9], [10], [13], [17]. The integral $\int_{0}^{\infty} f(s) \mathrm{d} Z_{s}$ is defined to be the limit in probability of $\int_{0}^{t} f(s) \mathrm{d} Z_{s}$ as $t \rightarrow \infty$ whenever it exists.

## 2. Characterization and properties of class $I(c)$

Throughout the paper fix $d$ in the set $\mathbb{N}$ of positive integers. Elements of $\mathbb{R}^{d}$ are column vectors. The inner product on $\mathbb{R}^{d}$ is denoted $\langle x, y\rangle$ and the corresponding norm is $|x|$. For a distribution $\mu$ on $\mathbb{R}^{d}$ denote the characteristic function of $\mu$ by $\widehat{\mu}$, $\widehat{\mu}(z)=\int_{\mathbb{R}^{d}} e^{\mathrm{i} i z, x\rangle} \mu(\mathrm{d} x), z \in \mathbb{R}^{d}$. For probability measures $\mu_{n}(n=1,2, \ldots)$ and $\mu$ on $\mathbb{R}^{d}, \mu_{n} \rightarrow \mu$ means weak convergence of $\mu_{n}$ to $\mu$. Let $D=\left\{x \in \mathbb{R}^{d}:|x| \leqslant 1\right\}$. For $\mu_{0} \in I D\left(\mathbb{R}^{d}\right)$ let $\Psi_{0}(z)=\log \widehat{\mu}_{0}(z)$, the distinguished logarithm of $\widehat{\mu}_{0}$ in [12], p. 33. We have

$$
\Psi_{0}(z)=-\frac{1}{2}\left\langle z, A_{0} z\right\rangle+\mathrm{i}\left\langle\gamma_{0}, z\right\rangle+\int_{\mathbb{R}^{d}}\left(e^{\mathrm{i}\langle z, x\rangle}-1-\mathrm{i}\langle z, x\rangle 1_{D}(x)\right) \nu_{0}(\mathrm{~d} x), \quad z \in \mathbb{R}^{d},
$$

where $\left(A_{0}, \nu_{0}, \gamma_{0}\right)$ is the triplet of $\mu_{0}$. For $s \geqslant 0$ let $\mu_{0}^{s}=\mu^{* s}$, which has characteristic function $\widehat{\mu}_{0}^{s}(z)=\exp s \Psi_{0}(z), z \in \mathbb{R}^{d}$.

If $f: \mathbb{R}_{+} \longrightarrow \mathbb{R}$ is bounded and measurable and $\left\{Z_{s}\right\}_{s \geqslant 0}$ is a Lévy process on $\mathbb{R}^{d}$, then the stochastic integral $\int_{a}^{b} f(s) \mathrm{d} Z_{s}$ is definable for $0 \leqslant a \leqslant b<\infty$ and has an infinitely divisible distribution on $\mathbb{R}^{d}$ with

$$
\begin{equation*}
E \exp \left(\mathrm{i}\left\langle z, \int_{a}^{b} f(s) \mathrm{d} Z_{s}\right\rangle\right)=\exp \int_{a}^{b} \Psi_{0}(f(s) z) \mathrm{d} s, \quad z \in \mathbb{R}^{d} \tag{2.1}
\end{equation*}
$$

Now and then we integrate with respect to a time-changed Lévy process. Hence, for $t>0$ and $s_{0} \geqslant 0, \int_{a}^{b} f(s) \mathrm{d}_{s} Z_{s_{0}+t s}$ denotes the stochastic integral of $f$ with respect to the Lévy process $\left\{Z_{s_{0}+t s}-Z_{s_{0}}\right\}_{s \geqslant 0}$.

For $m=1,2, \ldots$ let $I D_{\log ^{m}}\left(\mathbb{R}^{d}\right)$ be the class of distributions $\mu \in I D\left(\mathbb{R}^{d}\right)$ for which the Lévy measure $\nu$ of $\mu$ satisfies

$$
\begin{equation*}
\int_{|x|>2}(\log |x|)^{m} \nu(\mathrm{~d} x)<\infty \tag{2.2}
\end{equation*}
$$

Equivalently, $I D_{\log ^{m}}\left(\mathbb{R}^{d}\right)$ is the class of $\mu \in I D\left(\mathbb{R}^{d}\right)$ such that $\int_{|x|>2}(\log |x|)^{m} \mu(\mathrm{~d} x)<$ $\infty$. Let $I D_{\log }\left(\mathbb{R}^{d}\right)=I D_{\log ^{1}}\left(\mathbb{R}^{d}\right)$. The stochastic integral $\int_{0}^{\infty} e^{-c s} \mathrm{~d} Z_{s}$ with $c>0$ exists if and only if $\mathcal{L}\left(Z_{1}\right) \in I D_{\log }\left(\mathbb{R}^{d}\right)$. There is a one-to-one correspondence between $S D\left(\mathbb{R}^{d}\right)$ and $I D_{\log }\left(\mathbb{R}^{d}\right)$ as is mentioned in the previous section.

Let $\mu \in I D\left(\mathbb{R}^{d}\right)$ with triplet $(A, \nu, \gamma)$. Recall that if $\nu=0$ then $\mu$ is said to be Gaussian, and if $A=0$ then $\mu$ is said to be purely non-Gaussian. Let $M$ be a class of distributions on $\mathbb{R}^{d}$. Following [11] we say that $M$ is completely closed if the following three conditions (C1)-(C3) are satisfied, where (C1) $\mu_{1}, \mu_{2} \in M$ implies $\mu_{1} * \mu_{2} \in M$; (C2) $\mu_{n} \in M$ for $n=1,2, \ldots$ and $\mu_{n} \rightarrow \mu$ imply $\mu \in M$; (C3) $\mathcal{L}(X) \in M$ implies $\mathcal{L}(a X+b) \in M$ for all $a>0$ and $b \in \mathbb{R}^{d}$. Sometimes we also need the conditions (D), (P) and (GJ) on $M$. These are: (D) (dual property) if $\mathcal{L}(X) \in M$ then $\mathcal{L}(-X) \in M$; (P) (raising to the power) if $\mu \in M$ then $\mu^{t} \in M$ for all $t \geqslant 0$; (GJ) if $\mu \in M$ and $\mu=\mu_{G} * \mu_{J}$ where $\mu_{G}$ is Gaussian and $\mu_{J}$ is purely non-Gaussian, then $\mu_{G}, \mu_{J} \in M$. When we consider ( P ) or (GJ), we assume $M \subseteq I D\left(\mathbb{R}^{d}\right)$.

Definition 2.1. For $c \in \mathbb{R} \backslash\{0\}$ let $I(c)$ be the class of distributions given by

$$
\begin{equation*}
I(c)=\left\{\mathcal{L}\left(\int_{0}^{1} e^{c s} \mathrm{~d} Z_{s}\right):\left\{Z_{s}\right\}_{s \geqslant 0} \text { is a Lévy process on } \mathbb{R}^{d}\right\} \tag{2.3}
\end{equation*}
$$

We have $I(c) \subseteq I D\left(\mathbb{R}^{d}\right)$ as is mentioned above. We define the mapping $\Phi_{c}$ from $I D\left(\mathbb{R}^{d}\right)$ onto $I(c)$ by

$$
\begin{equation*}
\Phi_{c}\left(\mu_{0}\right)=\mathcal{L}\left(\int_{0}^{1} e^{c s} \mathrm{~d} Z_{s}\right), \quad \mu_{0}=\mathcal{L}\left(Z_{1}\right) \tag{2.4}
\end{equation*}
$$

Note that $\mu=\mathcal{L}\left(\int_{0}^{1} e^{c s} \mathrm{~d} Z_{s}\right) \in I(c)$ if and only if $\mu$ is the stationary distribution of a periodic Ornstein-Uhlenbeck process with parameter $c$ and background driving Lévy process $\left\{Z_{s}^{\prime}\right\}_{s \in[0,1]}$ where $Z_{s}^{\prime}=\left(e^{c}-1\right) Z_{s}$. See (1.4).

Remark 2.2. Let $t>0$ and $c \in \mathbb{R} \backslash\{0\}$. Then,

$$
I(c t)=\left\{\mathcal{L}\left(\int_{0}^{t} e^{c s} \mathrm{~d} Z_{s}\right):\left\{Z_{s}\right\}_{s \geqslant 0} \text { is a Lévy process on } \mathbb{R}^{d}\right\} .
$$

Indeed, $\int_{0}^{t} e^{c s} \mathrm{~d} Z_{s}=\int_{0}^{1} e^{c t s} \mathrm{~d}_{s} Z_{s t}$. As a special case,

$$
I(c)=\left\{\mathcal{L}\left(\int_{0}^{c} e^{s} \mathrm{~d} Z_{s}\right):\left\{Z_{s}\right\}_{s \geqslant 0} \text { is a Lévy process on } \mathbb{R}^{d}\right\}
$$

for $c>0$. Thus $I(c), c>0$, can be understood also as the class of stationary distributions of $c$-periodic Ornstein-Uhlenbeck processes with parameter 1 and background
driving Lévy process $\left\{Z_{s}^{\prime}\right\}_{s \in[0, c]}$, namely, the distributions of the solutions of

$$
\mathrm{d} X_{t}=-X_{t} \mathrm{~d} t+\mathrm{d} Z_{t}^{\prime}, \quad t \in[0, c], \quad \text { with } X_{0}=X_{c} \text { a.s. }
$$

Proposition 2.3. Let $c \in \mathbb{R} \backslash\{0\}$. Let $\mu_{0} \in I D\left(\mathbb{R}^{d}\right)$ have triplet $\left(A_{0}, \nu_{0}, \gamma_{0}\right)$ and corresponding distinguished logarithm $\Psi_{0}(z)$. Let $\mu=\Phi_{c}\left(\mu_{0}\right)$. Then,

$$
\begin{equation*}
\widehat{\mu}(z)=\exp \int_{0}^{1} \Psi_{0}\left(e^{c s} z\right) \mathrm{d} s, \quad z \in \mathbb{R}^{d} \tag{2.5}
\end{equation*}
$$

and the triplet of $\mu$ is $(A, \nu, \gamma)$, where

$$
\begin{align*}
A & =\int_{0}^{1} e^{2 c s} \mathrm{~d} s A_{0}=\frac{e^{2 c}-1}{2 c} A_{0},  \tag{2.6}\\
\nu(B) & =\int_{\mathbb{R}^{d}} \nu_{0}(\mathrm{~d} x) \int_{0}^{1} 1_{B}\left(e^{c s} x\right) \mathrm{d} s, \quad B \in \mathcal{B}\left(\mathbb{R}^{d}\right),  \tag{2.7}\\
\gamma & =\frac{e^{c}-1}{c} \gamma_{0}-\int_{\mathbb{R}^{d}} \nu_{0}(\mathrm{~d} x) \int_{0}^{1} e^{c s} x\left[1_{D}(x)-1_{D}\left(e^{c s} x\right)\right] \mathrm{d} s . \tag{2.8}
\end{align*}
$$

Proof. The expression for $\widehat{\mu}(z)$ follows from (2.1). The expression for the triplet is derived in [10], p. 35, or [13].

The following are basic properties of $I(c)$.
Proposition 2.4. Let $c \in \mathbb{R} \backslash\{0\}$.
(i) If $\mu \in I(c)$, then $\mu^{t} \in I(c)$ for all $t \geqslant 0$.
(ii) For $\mu_{0,1}, \mu_{0,2} \in I D\left(\mathbb{R}^{d}\right)$ we have $\Phi_{c}\left(\mu_{0,1}\right) * \Phi_{c}\left(\mu_{0,2}\right)=\Phi_{c}\left(\mu_{0,1} * \mu_{0,2}\right)$.
(iii) The class $I(c)$ satisfies (C1), (C3), (D), (P) and (GJ).
(iv) $I(c)=I(-c)$. In fact, if $\mu_{1}=\Phi_{c}\left(\mu_{0}\right)$ and $\mu_{2}=\Phi_{-c}\left(\mu_{0}\right)$, then $\widehat{\mu}_{2}(z)=$ $\widehat{\mu}_{1}\left(e^{-c} z\right), z \in \mathbb{R}^{d}$.

Proof. (i) Let $\mu=\mathcal{L}\left(\int_{0}^{1} e^{c s} \mathrm{~d} Z_{s}\right)$ and $\mu_{0}=\mathcal{L}\left(Z_{1}\right)$. Let $\Psi_{0}=\log \widehat{\mu}_{0}$. Then $\widehat{\mu}^{t}(z)=$ $\exp \left(t \int_{0}^{1} \Psi_{0}\left(e^{c s} z\right) \mathrm{d} s\right)$ by Proposition 2.3. It follows that $\mu^{t}=\mathcal{L}\left(\int_{0}^{1} e^{c s} \mathrm{~d}_{s} Z_{t s}\right) \in I(c)$. The proof of (ii) is left to the reader. By (i)-(ii) $I(c)$ satisfies (C1) and (P). It is, moreover, readily seen that $I(c)$ satisfies (C3), (D) and (GJ). The proof of (iv) is given by the use of (2.5).

Remark 2.5. As a special case of Proposition 2.10 (ii) below it follows that $I(c)$ satisfies (C2). Hence $I(c)$ is completely closed.

In the following we derive a characterization of $I(c)$. Let $S=\left\{\xi \in \mathbb{R}^{d}:|\xi|=1\right\}$ and Leb denote the Lebesgue measure on $(0, \infty)$.

Proposition 2.6 (Polar decomposition of a Lévy measure). Let $\nu$ be a Lévy measure on $\mathbb{R}^{d}$. Then, $\nu$ is decomposed in polar form as

$$
\nu(B)=\int_{S} \lambda(\mathrm{~d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \nu_{\xi}(\mathrm{d} u), \quad B \in \mathcal{B}\left(\mathbb{R}^{d}\right)
$$

where $\lambda$ is a finite measure on $S$ and $\nu_{\xi}$ is a measure on $(0, \infty)$ such that $\nu_{\xi}(B)$ is measurable in $\xi$ for each $B \in \mathcal{B}((0, \infty))$ and $\int_{0}^{\infty}\left(1 \wedge u^{2}\right) \nu_{\xi}(\mathrm{d} u)=1$ for $\xi \in \mathbb{R}^{d}$. This $\lambda$ is unique and $\nu_{\xi}$ is unique for $\lambda$-a.e. $\xi$.

Let $c>0$. If $\nu_{\xi}(\mathrm{d} u)$ has density $k_{\xi}(u) / u$ with respect to the Lebesgue measure on $(0, \infty)$ with some nonnegative $(\xi, u)$-measurable $k_{\xi}(u)$, then $\sum_{j=1}^{\infty} k_{\xi}\left(e^{j c} u\right)<\infty$ for $\lambda \times$ Leb-a.e. $(\xi, u)$.

Proof. The polar decomposition is constructed as in the proof of Theorem 15.10 in [12]. We leave the details to the reader.

Assume $\nu_{\xi}(\mathrm{d} u)$ has density $k_{\xi}(u) / u$. Denote $G_{\xi}^{*}(u)=\sum_{j=1}^{\infty} k_{\xi}\left(e^{j c} u\right)$. We have to show that $G_{\xi}^{*}(u)$ is finite for $\lambda \times$ Leb-a.e. $(\xi, u)$. For $a>0$ we have

$$
\begin{aligned}
& \int_{S} \lambda(\mathrm{~d} \xi) \int_{a}^{e^{c} a} \frac{G_{\xi}^{*}(u)}{u} \mathrm{~d} u=\sum_{j=1}^{\infty} \int_{S} \lambda(\mathrm{~d} \xi) \int_{a}^{e^{c} a} \frac{k_{\xi}\left(e^{j c} u\right)}{u} \mathrm{~d} u \\
= & \sum_{j=1}^{\infty} \int_{S} \lambda(\mathrm{~d} \xi) \int_{e^{j c} a}^{e^{(j+1) c} a} \frac{k_{\xi}(u)}{u} \mathrm{~d} u=\int_{S} \lambda(\mathrm{~d} \xi) \int_{e^{c} a}^{\infty} \frac{k_{\xi}(u)}{u} \mathrm{~d} u=\nu\left(\left\{x:|x|>e^{c a}\right\}\right)<\infty,
\end{aligned}
$$

which gives the desired result.

Theorem 2.7. Let $c>0$. Let $\mu \in I D\left(\mathbb{R}^{d}\right)$ have Lévy measure $\nu$ with polar decomposition $\lambda(\mathrm{d} \xi) \nu_{\xi}(\mathrm{d} u)$.
(i) $\mu$ is in $I(c)$ if and only if the following two conditions (a)-(b) are satisfied.
(a) For $\lambda$-a.e. $\xi$ the measure $\nu_{\xi}(\mathrm{d} u)$ has density $k_{\xi}(u) / u$ with respect to the Lebesgue measure on $(0, \infty)$ with some nonnegative $(\xi, u)$-measurable $k_{\xi}(u)$;
(b) Let $G_{\xi}^{*}(u)=\sum_{j=1}^{\infty} k_{\xi}\left(e^{j c} u\right)$. Then $G_{\xi}^{*}(u)$ has a version which is decreasing in $u$.
(ii) Assume $\mu \in I(c)$. Then there is one and only one $\mu_{0} \in I D\left(\mathbb{R}^{d}\right)$ such that $\mu=\Phi_{c}\left(\mu_{0}\right)$. There is a finite, right-continuous, decreasing version $G_{\xi}(u)$ of $G_{\xi}^{*}(u)$, and the relation between $\nu$ and the polar decomposition $\lambda_{0}(\mathrm{~d} \xi) \nu_{0 \xi}(\mathrm{~d} u)$
of the Lévy measure $\nu_{0}$ of $\mu_{0}$ is described as follows. We have

$$
\begin{gather*}
\nu_{0}(B)=-c \int_{S} \lambda(\mathrm{~d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \mathrm{d} G_{\xi}(u) \text { for } B \in \mathcal{B}\left(\mathbb{R}^{d}\right),  \tag{2.9}\\
k_{\xi}(u)=\frac{\nu_{0 \xi}\left(e^{-c} u, u\right]}{c a_{\xi}}, \quad G_{\xi}(u)=\frac{\nu_{0 \xi}(u, \infty)}{c a_{\xi}}, \quad \lambda(\mathrm{d} \xi)=a_{\xi} \lambda_{0}(\mathrm{~d} \xi), \tag{2.10}
\end{gather*}
$$

where

$$
\begin{equation*}
a_{\xi}=\int_{0}^{\infty}\left(1 \wedge u^{2}\right) \nu_{0 \xi}\left(e^{-c} u, u\right] \frac{\mathrm{d} u}{c u} . \tag{2.11}
\end{equation*}
$$

Proof. Let $(A, \nu, \gamma)$ be the triplet of $\mu$.
Step 1. Assume that $\nu$ satisfies (a)-(b). By Proposition 2.6 we have $G_{\xi}^{*}(u)<\infty$ almost surely. Therefore, $G_{\xi}^{*}(u)$ has a finite, decreasing, right-continuous version $G_{\xi}(u)$ by (b). For fixed $\xi,-G_{\xi}(u)$ induces a $\sigma$-finite measure on $(0, \infty)$. We show that the measure $\nu_{0}$ defined in (2.9) is a Lévy measure and that a distribution $\mu_{0} \in I D\left(\mathbb{R}^{d}\right)$ satisfies $\mu=\Phi_{c}\left(\mu_{0}\right)$ if it has triplet $\left(A_{0}, \nu_{0}, \gamma_{0}\right)$, where

$$
\begin{equation*}
A_{0}=\frac{2 c}{e^{2 c}-1} A, \quad \gamma_{0}=\frac{c}{e^{c}-1}\left(\gamma+\int_{\mathbb{R}^{d}} \nu_{0}(\mathrm{~d} x) \int_{0}^{1} e^{c s} x 1_{\left\{s: e^{-c s}<|x| \leqslant 1\right\}}(s) \mathrm{d} s\right) . \tag{2.12}
\end{equation*}
$$

We have (2.7). Indeed, $G_{\xi}\left(e^{-c} u\right)-G_{\xi}(u)=k_{\xi}(u)$ and hence with substitution $u=e^{c s} r$ we get

$$
\begin{aligned}
& \int_{\mathbb{R}^{d}} \nu_{0}(\mathrm{~d} x) \int_{0}^{1} 1_{B}\left(e^{c s} x\right) \mathrm{d} s=-c \int_{S} \lambda(\mathrm{~d} \xi) \int_{0}^{\infty} \mathrm{d} G_{\xi}(r) \int_{0}^{1} 1_{B}\left(e^{c s} r \xi\right) \mathrm{d} s \\
= & -c \int_{S} \lambda(\mathrm{~d} \xi) \int_{0}^{\infty} \mathrm{d} G_{\xi}(r) \int_{r}^{e^{c} r} 1_{B}(u \xi) \frac{\mathrm{d} u}{c u} \\
= & -\int_{S} \lambda(\mathrm{~d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \frac{\mathrm{d} u}{u} \int_{0}^{\infty} 1_{\left\{r: r \leqslant u<e^{c} r\right\}}(r) \mathrm{d} G_{\xi}(r) \\
= & \int_{S} \lambda(\mathrm{~d} \xi) \int_{0}^{\infty} 1_{B}(u \xi)\left(G_{\xi}\left(e^{-c} u\right)-G_{\xi}(u)\right) \frac{\mathrm{d} u}{u} \\
= & \int_{S} \lambda(\mathrm{~d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \frac{k_{\xi}(u)}{u} \mathrm{~d} u=\nu(B) .
\end{aligned}
$$

Notice that
$\int_{\mathbb{R}^{d}}\left(1 \wedge|x|^{2}\right) \nu_{0}(\mathrm{~d} x) \leqslant \int_{\mathbb{R}^{d}} \nu_{0}(\mathrm{~d} x) \int_{0}^{1}\left(1 \wedge\left|e^{c s} x\right|^{2}\right) \mathrm{d} s=\int_{\mathbb{R}^{d}}\left(1 \wedge|x|^{2}\right) \nu(\mathrm{d} x)<\infty$,
where the equality is by (2.7). Hence $\nu_{0}$ is a Lévy measure. Defining $A_{0}$ and $\gamma_{0}$ by (2.12) it is readily verified that we have (2.6) and (2.8). Hence the distribution $\mu_{0}$ with triplet $\left(A_{0}, \nu_{0}, \gamma_{0}\right)$ satisfies $\mu=\Phi_{c}\left(\mu_{0}\right)$.

Step 2. Let $\mu=\Phi_{c}\left(\mu_{0}\right)$, where $\mu_{0} \in I D\left(\mathbb{R}^{d}\right)$ has triplet $\left(A_{0}, \nu_{0}, \gamma_{0}\right)$. To complete the proof it suffices to show (a)-(b), (2.10)-(2.11), and that $\left(A_{0}, \nu_{0}, \gamma_{0}\right)$ satisfies (2.9)
and (2.12). (These imply the uniqueness of $\mu_{0}$, since $\lambda$ is unique and the function $G_{\xi}(u)$ is unique for $\lambda$-a.e. $\xi$.) Using (2.7) and substitution $e^{c s} r=u$, we get

$$
\begin{aligned}
\nu(B) & =\int_{\mathbb{R}^{d}} \nu_{0}(\mathrm{~d} x) \int_{0}^{1} 1_{B}\left(e^{c s} x\right) \mathrm{d} s=\int_{S} \lambda_{0}(\mathrm{~d} \xi) \int_{0}^{\infty} \nu_{0 \xi}(\mathrm{~d} r) \int_{0}^{1} 1_{B}\left(e^{c s} r \xi\right) \mathrm{d} s \\
& =\int_{S} \lambda_{0}(\mathrm{~d} \xi) \int_{0}^{\infty} \nu_{0 \xi}(\mathrm{~d} r) \int_{r}^{e^{c} r} 1_{B}(u \xi) \frac{\mathrm{d} u}{c u} \\
& =\int_{S} \lambda_{0}(\mathrm{~d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \frac{\nu_{0 \xi}\left(e^{-c} u, u\right]}{c u} \mathrm{~d} u .
\end{aligned}
$$

Hence we have (a). By the normalization $\int_{0}^{\infty}\left(1 \wedge r^{2}\right) \nu_{\xi}(\mathrm{d} r)=1$ we get the representations of $\lambda(\mathrm{d} \xi)$ and $k_{\xi}(u)$ in (2.10)-(2.11). Since $k_{\xi}\left(e^{j c} u\right)=\left(c a_{\xi}\right)^{-1} \nu_{0 \xi}\left(e^{(j-1) c} u, e^{j c} u\right]$ we have $\sum_{j=1}^{\infty} k_{\xi}\left(e^{j c} u\right)=\left(c a_{\xi}\right)^{-1} \nu_{0 \xi}(u, \infty)$. Obviously $G_{\xi}(u):=\left(c a_{\xi}\right)^{-1} \nu_{0 \xi}(u, \infty)$ is finite, right-continuous and decreasing. In particular we have (b). The expression (2.12) follows from (2.6) and (2.8). Finally, in order to show (2.9) note that we have

$$
\begin{aligned}
& -c \int_{S} \lambda(\mathrm{~d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \mathrm{d} G_{\xi}(u)=c \int_{S} \lambda(\mathrm{~d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \frac{\nu_{0 \xi}(\mathrm{~d} u)}{c a_{\xi}} \\
= & \int_{S} \lambda_{0}(\mathrm{~d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \nu_{0 \xi}(\mathrm{~d} u)=\nu_{0}(B)
\end{aligned}
$$

by (2.10)-(2.11).
Corollary 2.8. For $c \in \mathbb{R} \backslash\{0\}$ the mapping $\Phi_{c}: I D\left(\mathbb{R}^{d}\right) \longrightarrow I(c)$ is one-to-one.
To prove the corollary note that we may assume $c>0$ by Proposition 2.4 (iv), in which case the result follows from (ii) above. This corollary shows that if $\left\{X_{s}^{j}\right\}_{s \in[0,1]}$, $j=1,2$, are periodic Ornstein-Uhlenbeck processes with parameter $c>0$ and background driving Lévy processes $\left\{Z_{s}^{j}\right\}_{s \in[0,1]}$ such that $\mathcal{L}\left(X_{s}^{1}\right)=\mathcal{L}\left(X_{s}^{2}\right)$, then $\mathcal{L}\left(Z_{1}^{1}\right)=$ $\mathcal{L}\left(Z_{1}^{2}\right)$. That is, the stationary distribution of a periodic Ornstein-Uhlenbeck process determines the law of the background driving Lévy process uniquely.

Remark 2.9. In the next result it is convenient to let $I(0)=I D\left(\mathbb{R}^{d}\right)$. Also let $\Phi_{0}: I(0) \longrightarrow I(0)$ be the identity mapping.

Proposition 2.10. Let $c_{n}$ and $c$ be real numbers with $c_{n} \rightarrow c$.
(i) Let $\mu_{n, 0} \in I D\left(\mathbb{R}^{d}\right)$ for $n=1,2, \ldots$ and $\mu_{n, 0} \rightarrow \mu_{0}$. Then $\mu_{0} \in I D\left(\mathbb{R}^{d}\right)$ and $\Phi_{c_{n}}\left(\mu_{n, 0}\right) \rightarrow \Phi_{c}\left(\mu_{0}\right)$.
(ii) Conversely, let $\mu_{n} \in I\left(c_{n}\right)$ and $\mu_{n} \rightarrow \mu$. Then $\mu \in I(c)$ and $\Phi_{c_{n}}^{-1}\left(\mu_{n}\right) \rightarrow$ $\Phi_{c}^{-1}(\mu)$.

Before proving this result we need two lemmas.

Lemma 2.11. Let $c>0$ and $\mu=\Phi_{c}\left(\mu_{0}\right)$, where $\mu$ and $\mu_{0}$ have triplets $\left(A_{,} \nu, \gamma\right)$ and $\left(A_{0}, \nu_{0}, \gamma_{0}\right)$, respectively. Then we have

$$
\begin{align*}
& \int_{|x|>a} \nu_{0}(\mathrm{~d} x) \leqslant \int_{|x|>a} \nu(\mathrm{~d} x) \leqslant \int_{|x|>e^{-c_{a}}} \nu_{0}(\mathrm{~d} x), \quad a>0,  \tag{2.13}\\
& \frac{e^{2 c}-1}{2 c} \int_{|x| \leqslant e^{-c}}|x|^{2} \nu_{0}(\mathrm{~d} x) \leqslant \int_{|x| \leqslant 1}|x|^{2} \nu(\mathrm{~d} x) \leqslant \frac{e^{2 c}-1}{2 c} \int_{|x| \leqslant 1}|x|^{2} \nu_{0}(\mathrm{~d} x),  \tag{2.14}\\
& \left|\int_{\mathbb{R}^{d}} \nu_{0}(\mathrm{~d} x) \int_{0}^{1} e^{c s} x 1_{\left\{x: e^{-c s}<|x| \leqslant 1\right\}}(x) \mathrm{d} s\right| \leqslant \frac{e^{c}-1}{c} \int_{e^{-c}<|x| \leqslant 1} \nu_{0}(\mathrm{~d} x) . \tag{2.15}
\end{align*}
$$

Proof. For $a>0$ let $D^{a}=\left\{y \in \mathbb{R}^{d}:|y|>a\right\}$. Then, for $s \in[0,1]$ we have $1_{D^{a}}(x) \leqslant$ $1_{D^{a}}\left(e^{c s} x\right) \leqslant 1_{D^{a}}\left(e^{c} x\right)$ and hence

$$
\int_{D^{a}} \nu_{0}(\mathrm{~d} x) \leqslant \int_{\mathbb{R}^{d}} \nu_{0}(\mathrm{~d} x) \int_{0}^{1} 1_{D^{a}}\left(e^{c s} x\right) \mathrm{d} s \leqslant \int_{\mathbb{R}^{d}} 1_{D^{a}}\left(e^{c} x\right) \nu_{0}(\mathrm{~d} x)
$$

from which (2.13) follows since $\int_{\mathbb{R}^{d}} \nu_{0}(\mathrm{~d} x) \int_{0}^{1} 1_{D^{a}}\left(e^{c s} x\right) \mathrm{d} s=\nu\left(D^{a}\right)$ by (2.7). Note that we have

$$
\begin{equation*}
\int_{\mathbb{R}^{d}} f(x) \nu(\mathrm{d} x)=\int_{\mathbb{R}^{d}} \nu_{0}(\mathrm{~d} x) \int_{0}^{1} f\left(e^{c s} x\right) \mathrm{d} s \tag{2.16}
\end{equation*}
$$

for all nonnegative measurable functions $f$ by (2.7). Thus (2.14) follows. Since $|x| 1_{\left\{e^{-c s}<|x| \leqslant 1\right\}}(x) \leqslant 1_{\left\{e^{-c<|x| \leqslant 1\}}\right.}(x)$ for $0<s<1$, (2.15) is obvious.

Remark 2.12. Let $c>0$ and let $\mu, \mu_{0}, \nu, \nu_{0}$ be as in Lemma 2.11. We can also prove the following from (2.16).
(i) $\nu\left(\mathbb{R}^{d}\right)<\infty$ if and only if $\nu_{0}\left(\mathbb{R}^{d}\right)<\infty$.
(ii) Let $0<\alpha<2$. Then $\int_{|x| \leqslant 1}|x|^{\alpha} \nu(\mathrm{d} x)<\infty$ if and only if $\int_{|x| \leqslant 1}|x|^{\alpha} \nu_{0}(\mathrm{~d} x)<\infty$.
(iii) Let $\alpha>0$. Then $\int_{|x|>1}|x|^{\alpha} \nu(\mathrm{d} x)<\infty$ if and only if $\int_{|x|>1}|x|^{\alpha} \nu_{0}(\mathrm{~d} x)<\infty$.
(iv) Let $\alpha>0$. Then $\int_{|x|>1}(\log |x|)^{\alpha} \nu(\mathrm{d} x)<\infty$ if and only if $\int_{|x|>1}(\log |x|)^{\alpha} \nu_{0}(\mathrm{~d} x)<$ $\infty$.

Recall that a set of probability measures on $\mathbb{R}^{d}$ is precompact if and only if any sequence in the set has a subsequence which converges weakly.

Lemma 2.13. Let $M$ be a subclass of $I D\left(\mathbb{R}^{d}\right)$ and let $C$ be a bounded set in $\mathbb{R}$. If $M$ is precompact then $\left\{\Phi_{c}\left(\mu_{0}\right): \mu_{0} \in M, c \in C\right\}$ and $\left\{\Phi_{c}^{-1}(\mu): \mu \in M, c \in C\right\}$ are both precompact.

Proof. Let $N$ denote a subset of $I D\left(\mathbb{R}^{d}\right)$ and $\mu \in N$ have triplet $\left(A_{\mu}, \nu_{\mu}, \gamma_{\mu}\right)$. Then, $N$ is precompact if and only if the following four conditions (i)-(iv) are satisfied, where (i) $\sup _{\mu \in N} \sup _{|x| \leqslant 1}\left|A_{\mu} x\right|<\infty$; (ii) $\sup _{\mu \in N} \int_{\mathbb{R}^{d}}\left(|x|^{2} \wedge 1\right) \nu_{\mu}(\mathrm{d} x)<\infty$; (iii) $\sup _{\mu \in N}\left|\gamma_{\mu}\right|<\infty$;
(iv) $\lim _{l \rightarrow \infty} \sup _{\mu \in N} \nu_{\mu}(\{|x|>l\})=0$. Indeed, this is essentially E 12.5 in [12] except that (iv) is missing in that exercise. Hence, we have to show that (i)-(iv) are satisfied with $N=\left\{\Phi_{c}\left(\mu_{0}\right): \mu_{0} \in M, c \in C\right\}$ and $\left\{\Phi_{c}^{-1}(\mu): \mu \in M, c \in C\right\}$ if (i)-(iv) are satisfied with $N=M$. This is, however, easily verified using (2.6)-(2.8) and Lemma 2.11 if $C \subseteq[0, \infty)$. Further, use Proposition 2.4 (iv) if $C \cap(-\infty, 0) \neq \emptyset$.

Proof of Proposition 2.10. (i) By [12], Lemma 7.8, $\mu_{0} \in I D\left(\mathbb{R}^{d}\right)$. Let $\Psi_{0}(z)=$ $\log \widehat{\mu}_{0}(z)$ and $\Psi_{n, 0}(z)=\log \widehat{\mu_{n, 0}}(z)$. By Lemma 7.7 in [12], $\Psi_{n, 0}(z) \rightarrow \Psi_{0}(z)$ uniformly on any compact set. Hence $\int_{0}^{1} \Psi_{n, 0}\left(e^{c_{n} s} z\right) \mathrm{d} s \rightarrow \int_{0}^{1} \Psi_{0}\left(e^{c s} z\right) \mathrm{d} s$ for all $z \in \mathbb{R}^{d}$. Hence, by $(2.5) \widehat{\Phi_{c_{n}}\left(\mu_{n, 0}\right)}(z) \rightarrow \widehat{\Phi_{c}\left(\mu_{0}\right)}(z)$. Note that $(2.5)$ is trivially true for $c=0$.
(ii) $\left\{\mu_{n}: n=1,2, \ldots\right\}$ is precompact. Hence $\left\{\Phi_{c_{n}}^{-1}\left(\mu_{n}\right): n=1,2, \ldots\right\}$ is precompact by Lemma 2.13. Thus, there is a subsequence $\left\{\Phi_{c_{n^{\prime}}}^{-1}\left(\mu_{n^{\prime}}\right)\right\}$ such that $\Phi_{c_{n^{\prime}}}^{-1}\left(\mu_{n^{\prime}}\right) \rightarrow$ $\mu_{0}$ where $\mu_{0}$ is a probability measure on $\mathbb{R}^{d}$. We have $\mu_{0} \in I D\left(\mathbb{R}^{d}\right)$ and by (i) $\mu_{n^{\prime}} \rightarrow \Phi_{c}\left(\mu_{0}\right)$. Hence $\mu_{0}=\Phi_{c}^{-1}(\mu)$. Thus $\mu_{0}$ is independent of the choice of subsequence. Hence $\Phi_{c_{n}}^{-1}\left(\mu_{n}\right) \rightarrow \Phi_{c}^{-1}(\mu)$.

## 3. Classes $I(c)$ for varying values of $c$

In this section we study the relations between classes $I(c)$ for different values of $c$. First we compare $I(c)$ to $S D\left(\mathbb{R}^{d}\right)$.

Theorem 3.1. (i) We have $S D\left(\mathbb{R}^{d}\right) \varsubsetneqq I(c)$ for any $c \in \mathbb{R} \backslash\{0\}$.
(ii) Let $c_{n} \in \mathbb{R}$ with $c_{n} \uparrow \infty$. Then $\bigcap_{n} I\left(c_{n}\right)=S D\left(\mathbb{R}^{d}\right)$.

Remark 3.2. Recall that a distribution is selfdecomposable if and only if it is the stationary distribution of a usual Ornstein-Uhlenbeck process (driven by a Lévy process). Hence the result in (i) shows that a stationary distribution of a usual OrnsteinUhlenbeck process is, for any $c \neq 0$, the stationary distribution of a periodic OrnsteinUhlenbeck process with parameter $c$ as well.

Proof of Theorem 3.1. (i) We may and do assume $c>0$ without loss of generality. Let $\mu \in S D\left(\mathbb{R}^{d}\right)$. Then $\mu$ satisfies condition (a) of Theorem 2.7 and $k_{\xi}(u)$ is decreasing in $u$, see [12], Theorem 15.10. Hence also (b) of Theorem 2.7 is satisfied, which means $\mu \in I(c)$. If $\mu \in S D\left(\mathbb{R}^{d}\right)$ its Lévy measure has total mass 0 or $\infty$. Hence Remark 2.12 (i) implies $S D\left(\mathbb{R}^{d}\right) \varsubsetneqq I(c)$.
(ii) Step 1. Let $\mu=\mathcal{L}\left(\int_{0}^{c} e^{-s} \mathrm{~d} Z_{s}^{(c)}\right)$ with $0<c<\infty$, where $\left\{Z_{s}^{(c)}\right\}_{s \geqslant 0}$ is a Lévy process on $\mathbb{R}^{d}$. Then the Lévy measures $\nu$ and $\nu^{(c)}$ of $\mu$ and $\mathcal{L}\left(Z_{1}^{(c)}\right)$ are related as

$$
\begin{equation*}
\nu(B)=\int_{\mathbb{R}^{d}} \nu^{(c)}(\mathrm{d} x) \int_{0}^{c} 1_{B}\left(e^{-s} x\right) \mathrm{d} s, \quad B \in \mathcal{B}\left(\mathbb{R}^{d}\right) \tag{3.1}
\end{equation*}
$$

Let $\lambda^{(c)}(\mathrm{d} \xi) \nu_{\xi}^{(c)}(\mathrm{d} r)$ be the polar decomposition of $\nu^{(c)}$. As in the second step of the proof of Theorem 2.7 one shows that

$$
\begin{equation*}
\nu(B)=\int_{S} \lambda^{(c)}(\mathrm{d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \nu_{\xi}^{(c)}\left(u, e^{c} u\right] \frac{\mathrm{d} u}{u} . \tag{3.2}
\end{equation*}
$$

Notice that if $B \subseteq\{x: \alpha<|x| \leq \beta\}$ with $0<\alpha<\beta<\infty$, then

$$
\begin{aligned}
& \int_{S} \lambda^{(c)}(\mathrm{d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \nu_{\xi}^{(c)}(u, \infty) \frac{\mathrm{d} u}{u} \leqslant \int_{S} \lambda^{(c)}(\mathrm{d} \xi) \int_{\alpha}^{\beta} \nu_{\xi}^{(c)}(u, \infty) \frac{\mathrm{d} u}{u} \\
\leqslant & \frac{\beta-\alpha}{\alpha} \int_{S} \lambda^{(c)}(\mathrm{d} \xi) \nu_{\xi}^{(c)}(\alpha, \infty)<\infty
\end{aligned}
$$

It follows from (3.2) that

$$
\begin{equation*}
\nu(B)=\int_{S} \lambda^{(c)}(\mathrm{d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \nu_{\xi}^{(c)}(u, \infty) \frac{\mathrm{d} u}{u}-\int_{S} \lambda^{(c)}(\mathrm{d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \nu_{\xi}^{(c)}\left(e^{c} u, \infty\right) \frac{\mathrm{d} u}{u} \tag{3.3}
\end{equation*}
$$

if $B \subseteq\{x: \alpha<|x| \leqslant \beta\}$ for some $0<\alpha<\beta<\infty$.
Step 2. Let $c, \nu$ and $\nu^{(c)}$ be as in Step 1. For any $l>0$ we have $\int_{|x|>l} \nu(\mathrm{~d} x) \geqslant$ $c \int_{|x|>e^{c} l} \nu^{(c)}(\mathrm{d} x)$. Indeed, using (3.2) and $D^{l}=\{y:|y|>l\}$, we get

$$
\begin{aligned}
\int_{|x|>l} \nu(\mathrm{~d} x) & =\int_{\mathbb{R}^{d}} \nu^{(c)}(\mathrm{d} x) \int_{0}^{c} 1_{D^{l}}\left(e^{-s} x\right) \mathrm{d} s=c \int_{\mathbb{R}^{d}} \nu^{(c)}(\mathrm{d} x) \int_{0}^{1} 1_{D^{l}}\left(e^{-c s} x\right) \mathrm{d} s \\
& =c \int_{\mathbb{R}^{d}} \nu^{(c)}(\mathrm{d} x) \int_{0}^{1} 1_{D^{l e^{c}}}\left(e^{c s} x\right) \mathrm{d} s \geqslant c \int_{\mathbb{R}^{d}} 1_{D^{l e c}}(x) \nu^{(c)}(\mathrm{d} x) .
\end{aligned}
$$

Step 3. Assume that $\mu \in \bigcap_{n=1}^{\infty} I\left(c_{n}\right)$ for some $c_{n} \uparrow \infty$. Since $I\left(c_{n}\right)=I\left(-c_{n}\right)$, $\mu$ has the representation $\mu=\mathcal{L}\left(\int_{0}^{c_{n}} e^{-s} \mathrm{~d} Z_{s}^{\left(c_{n}\right)}\right)$ by Remark 2.2. Let $B \in \mathcal{B}\left(\mathbb{R}^{d}\right)$ satisfying $B \subseteq\{x: \alpha<|x| \leqslant \beta\}$ with $0<\alpha<\beta<\infty$. By Step 2 we see that

$$
\begin{aligned}
& \int_{S} \lambda^{\left(c_{n}\right)}(\mathrm{d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \nu_{\xi}^{\left(c_{n}\right)}\left(e^{c_{n}} u, \infty\right) \frac{\mathrm{d} u}{u} \leqslant \int_{S} \lambda^{\left(c_{n}\right)}(\mathrm{d} \xi) \int_{\alpha}^{\beta} \nu_{\xi}^{\left(c_{n}\right)}\left(e^{c_{n}} u, \infty\right) \frac{\mathrm{d} u}{u} \\
\leqslant & \frac{\beta-\alpha}{\alpha} \int_{S} \lambda^{\left(c_{n}\right)}(\mathrm{d} \xi) \nu_{\xi}^{\left(c_{n}\right)}\left(e^{c_{n}} \alpha, \infty\right)=\frac{\beta-\alpha}{\alpha} \nu^{\left(c_{n}\right)}\left(\left\{x:|x|>e^{c_{n}} \alpha\right\}\right) \\
\leqslant & \frac{\beta-\alpha}{\alpha c_{n}} \nu(\{x:|x|>\alpha\}) \rightarrow 0 \quad \text { as } n \rightarrow \infty .
\end{aligned}
$$

Hence, by (3.3),

$$
\begin{equation*}
\nu(B)=\lim _{n \rightarrow \infty} \int_{S} \lambda^{\left(c_{n}\right)}(\mathrm{d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \nu_{\xi}^{\left(c_{n}\right)}(u, \infty) \frac{\mathrm{d} u}{u} \tag{3.4}
\end{equation*}
$$

It follows that, for $b>1$,

$$
\begin{aligned}
\nu\left(b^{-1} B\right) & =\lim _{n \rightarrow \infty} \int_{S} \lambda^{\left(c_{n}\right)}(\mathrm{d} \xi) \int_{0}^{\infty} 1_{B}(b u \xi) \nu_{\xi}^{\left(c_{n}\right)}(u, \infty) \frac{\mathrm{d} u}{u} \\
& =\lim _{n \rightarrow \infty} \int_{S} \lambda^{\left(c_{n}\right)}(\mathrm{d} \xi) \int_{0}^{\infty} 1_{B}(u \xi) \nu_{\xi}^{\left(c_{n}\right)}\left(b^{-1} u, \infty\right) \frac{\mathrm{d} u}{u}
\end{aligned}
$$

Since $\nu_{\xi}^{\left(c_{n}\right)}\left(b^{-1} u, \infty\right) \geqslant \nu_{\xi}^{\left(c_{n}\right)}(u, \infty)$, we get $\nu\left(b^{-1} B\right) \geqslant \nu(B)$ for every $B \in \mathcal{B}\left(\mathbb{R}^{d}\right)$ satisfying $B \subseteq\{x: \alpha<|x| \leqslant \beta\}$ with $0<\alpha<\beta<\infty$. This shows that $\mu$ is selfdecomposable by [12], Theorem 15.8.

For $\mu \in I D\left(\mathbb{R}^{d}\right)$ let us introduce $Q(\mu)$, the set of $c \geqslant 0$ for which $\mu$ is in $I(c)$. That is, $Q(\mu)=\{c \geqslant 0: \mu \in I(c)\}$. Here we let $I(0)=I D\left(\mathbb{R}^{d}\right)$ as in Remark 2.9. Hence $0 \in Q(\mu)$. Due to Proposition 2.4 (iv) we consider only nonnegative values of c.

Proposition 3.3. Let $\mu \in I D\left(\mathbb{R}^{d}\right)$.
(i) The set $Q(\mu)$ is a closed subset of $\mathbb{R}_{+}$;
(ii) if $c \in Q(\mu)$ then $c / n \in Q(\mu)$ for all $n \in \mathbb{N}$;
(iii) the set $Q(\mu)$ is either equal to $[0, \infty)$ or bounded;
(iv) we have $Q(\mu)=[0, \infty)$ if and only if $\mu \in S D\left(\mathbb{R}^{d}\right)$.

Proof. The fact (i) follows from Proposition 2.10 while (iii) and (iv) are due to Theorem 3.1; (ii) is equivalent to the assertion that $I(c) \subseteq I(c / n)$ for $n \in \mathbb{N}$. Assume that $c>0$ and $\mu \in I(c)$. Let us prove that $\mu \in I\left(c^{\prime}\right)$ for $c^{\prime}=c / n$. Denote $G_{\xi, l}^{*}(u)=\sum_{j=1}^{\infty} k_{\xi}\left(e^{(j n-l) c^{\prime}} u\right)$ for $l=0,1, \ldots, n-1$, where $k_{\xi}$ is defined in Theorem 2.7. By Theorem $2.7 G_{\xi, l}^{*}(u)$ has a decreasing version. Since $G_{\xi, 0}^{*}(u)+\cdots+G_{\xi, n-1}^{*}(u)=$ $\sum_{j=1}^{\infty} k_{\xi}\left(e^{j c^{\prime}} u\right)$ we have $\mu \in I\left(c^{\prime}\right)$, using Theorem 2.7 again.

Let us consider some special distributions to show the rich variety of the set $Q(\mu)$.
Proposition 3.4. Let $d=1$.
(i) Let $\mu$ be an infinitely divisible distribution $\mathbb{R}$ for which the Lévy measure is not absolutely continuous with respect to the Lebesgue measure on $\mathbb{R}$. Then $Q(\mu)=\{0\}$.

In the following let $\mu \in I D(\mathbb{R})$ have Lévy measure $\nu$ given by $\nu(\mathrm{d} u)=1_{(0, \infty)}(u) \frac{k(u)}{u} \mathrm{~d} u$. Let $l(v)=k\left(e^{v}\right)$ for $v \in \mathbb{R}$.
(ii) If $c>0$ and $\mu=\Phi_{c}\left(\mu_{0}\right)$, where $\mu_{0}$ is the Poisson distribution with mean 1 , then $l(v)=\frac{1}{c} 1_{[0, c)}(v)$ and $Q(\mu)=\{0\} \cup\{c / n: n=1,2, \ldots\}$.
(iii) Assume $l(v)$ is either strictly increasing or strictly decreasing on an interval $[a, b)$ with $a<b$ and that $l(v)=0$ for $v \in \mathbb{R} \backslash[a, b)$. Then $Q(\mu)=\{0\}$.
(iv) Let $l(v)$ be decreasing on $(-\infty, 0]$ with $l(0)=2$ and

$$
l(v)= \begin{cases}v+2, & v \in[0,1]  \tag{3.5}\\ 3-3(v-1), & v \in[1,2] \\ 0, & v \geqslant 2\end{cases}
$$

Then $Q(\mu)=[0,1]$.
(v) Let $l(v)$ be decreasing on $(-\infty, 0]$ with $l(0)=4$ and

$$
l(v)= \begin{cases}4+2 v, & v \in[0,1]  \tag{3.6}\\ 6-6(v-1), & v \in[1,3 / 2] \\ 3-3(v-3 / 2), & v \in[3 / 2,5 / 2] \\ 0, & v \geqslant 5 / 2\end{cases}
$$

Then $Q(\mu)=[0,3 / 4] \cup[1,3 / 2]$.

Proof. (i) For $c>0$ a necessary condition for a distribution to be in $I(c)$ is by Theorem 2.7 that its Lévy measure is absolutely continuous.

To prove (ii)-(v) let $M_{c}(v)=\sum_{j=0}^{\infty} l(v+j c)$ for $c>0$ and $v \in \mathbb{R}$.
(ii) The Lévy measure of $\mu_{0}$ is $\nu_{0}=\delta_{1}$, which by Theorem 2.7 implies that the Lévy measure of $\mu$ is $\nu(\mathrm{d} u)=1_{(0, \infty)}(u) \frac{k(u)}{u} \mathrm{~d} u$, where $k(u)=\frac{1}{c} 1_{\left[1, e^{c}\right)}(u)$. Hence $l(v)=$ $\frac{1}{c} 1_{[0, c)}(v)$. From Proposition 3.3 (ii) it follows that $Q(\mu)$ includes $\{0\} \cup\{c / n: n=$ $1,2, \ldots\}$. Let $c^{\prime} \in(0, \infty) \backslash\{c / n: n=1,2, \ldots\}$. In order to show $c^{\prime} \notin Q(\mu)$ it suffices by Theorem 2.7 to show that $M_{c^{\prime}}(v)$ does not have a decreasing version on $\mathbb{R}$. If $c^{\prime}>c$, then $M_{c^{\prime}}(v)$ is 0 on $\left[c-c^{\prime}, 0\right)$ and $1 / c$ on $[0, c)$, which means $c^{\prime} \notin Q(\mu)$. If $0<c^{\prime}<c$ with $c^{\prime} \notin\{c / n: n=1,2, \ldots\}$, then choosing $n \in \mathbb{N}$ such that $c /(n+1)<c^{\prime}<c / n$, we see that $M_{c^{\prime}}(v)$ equals $n / c$ on $\left(c-(n+1) c^{\prime}, 0\right)$ and $(n+1) / c$ on $\left(0, c-n c^{\prime}\right)$, which means $c^{\prime} \notin Q(\mu)$.
(iii) First we consider the case where $l(v)$ is strictly increasing on $[a, b)$. Let $c>0$. By Theorem 2.7 we have to show that $M_{c}(v)$ does not have a decreasing version. If $a \vee(b-c)<v_{1}<v_{2}<b$ then $M_{c}\left(v_{1}\right)=l\left(v_{1}\right)<l\left(v_{2}\right)=M\left(v_{2}\right)$ and the result follows.

Next assume that $l(v)$ is strictly decreasing on $[a, b)$. It is readily seen that $M_{c}(v)$ is strictly decreasing on $[a, b)$. For $v<a$ we have

$$
M_{c}(v)=\sum_{j=0}^{\infty} l(v+j c)=\sum_{j=1}^{\infty} l(v+j c)=M_{c}(v+c) .
$$

Let $c>0$. If $v_{1}<a<v_{2}<v_{1}+c<b$, then $M_{c}\left(v_{1}\right)=M_{c}\left(v_{1}+c\right)<M_{c}\left(v_{2}\right)$. Hence $c \notin Q(\mu)$.
(iv) It suffices to show that
(a) $v \mapsto M_{c}(v)$ is decreasing for $c \in[1 / 2,1]$,
(b) $v \mapsto M_{c}(v)$ does not have a decreasing version for $c>1$.

Indeed, by Theorem 2.7 and Proposition 3.3 (ii), (a) and (b) imply that $c \notin Q(\mu)$ for $c>1$ and $c / n \in Q(\mu)$ for $c \in[1 / 2,1]$ and $n \in \mathbb{N}$. It follows that $Q(\mu)=[0,1]$.

Proof of (a). Let $c \in[1 / 2,1]$. We consider four cases.
Case 1: $v \in[1, \infty)$. Note that $v \mapsto l(v+j c)$ is decreasing on $[1, \infty)$ for $j \geqslant 0$. Thus $M_{c}(v)$ is decreasing in $v$ for $v \in[1, \infty)$.
Case 2: $v \in[1-c, 1]$. Similar to case 3 below and hence omitted.
Case 3: $v \in[0,1-c]$. Let $0 \leqslant v_{1} \leqslant v_{2} \leqslant 1-c$. Then $0 \leqslant v_{1}+c \leqslant v_{2}+c \leqslant$ $1,1 \leqslant v_{1}+2 c \leqslant v_{2}+2 c \leqslant 2$ and $1 \leqslant v_{1}+j c \leqslant v_{2}+j c$ for $j \geqslant 2$. Hence $l\left(v_{2}\right)-l\left(v_{1}\right)=l\left(v_{2}+c\right)-l\left(v_{1}+c\right)=v_{2}-v_{1}, l\left(v_{2}+2 c\right)-l\left(v_{1}+2 c\right)=-3\left(v_{2}-v_{1}\right)$ and $l\left(v_{2}+j c\right)-l\left(v_{1}+j c\right) \leqslant 0$ for $j \geqslant 2$. Therefore $M_{c}\left(v_{2}\right)-M_{c}\left(v_{1}\right) \leqslant-\left(v_{2}-v_{1}\right) \leqslant 0$.
Case 4: $v \in(-\infty, 0]$. Let $v_{1}, v_{2}$ satisfy $-c \leqslant v_{1} \leqslant v_{2} \leqslant 0$. Then $l\left(v_{1}\right) \geqslant l\left(v_{2}\right)$ since $l(v)$ is decreasing on $(-\infty, 0]$, and $M_{c}\left(v_{1}+c\right) \geqslant M_{c}\left(v_{2}+c\right)$ since $0 \leqslant v_{1}+c \leqslant v_{2}+c$ and $M_{c}(v)$ is decreasing on $[0, \infty)$. Since $M_{c}(v)=l(v)+M_{c}(v+c)$, we have $M_{c}\left(v_{1}\right) \geqslant$ $M_{c}\left(v_{2}\right)$, which shows that $M_{c}(v)$ is decreasing on $[-c, 0]$. Repeating inductively it follows that $M_{c}(v)$ is decreasing on any interval $[-c p,-c(p-1)]$ for $p \in \mathbb{N}$.

Proof of (b). Let $c>1$. If $2-c<v_{1}<v_{2}<1$, then $M_{c}\left(v_{1}\right)=l\left(v_{1}\right)<l\left(v_{2}\right)=$ $M_{c}\left(v_{2}\right)$.
(v) We show the following.
(c) $v \mapsto M_{c}(v)$ is decreasing for $c \in[1 / 4,1 / 2]$.
(d) $v \mapsto M_{c}(v)$ does not have a decreasing version for $c \in(3 / 4,1)$.
(e) $v \mapsto M_{c}(v)$ is decreasing for $c \in[1,3 / 2]$.
(f) $v \mapsto M_{c}(v)$ does not have a decreasing version for $c>3 / 2$.

Since $c \in Q(\mu)$ implies $c / n \in Q(\mu)$ for all $n \in \mathbb{N}$, it follows from (c) and (e) that $[0,1 / 2]$ and $[1 / 2,3 / 4]$ are subsets of $Q(\mu)$. Hence, it follows from (c)-(f) that $Q(\mu)=$ $[0,3 / 4] \cup[1,3 / 2]$.

Proof of (c). Fix $c \in[1 / 4,1 / 2]$. Clearly $M_{c}(v)$ is decreasing on $[1, \infty)$. Let $q_{0} \in\{2,3,4\}$ be the least $q$ such that $q c \geqslant 1$. We prove that $M_{c}(v)$ is decreasing on the interval $[(1-q c) \vee 0,1-(q-1) c]$ for $q=1, \ldots, q_{0}$. Since $\left(1-q_{0} c\right) \vee 0=0$, this implies that $M_{c}(v)$ is decreasing on $[0,1]$, and (c) follows from decreasingness of $l(v)$ on $(-\infty, 0]$. Hence let $q \in\left\{1, \ldots, q_{0}\right\}$ and $v_{1}, v_{2} \in[(1-q c) \vee 0,1-(q-1) c]$ with $v_{1} \leqslant v_{2}$. For $p=0, \ldots, q-1$ and $i=1,2$ we have $v_{i}+p c \in[0,1]$. Hence $l\left(v_{1}+p c\right)-l\left(v_{2}+p c\right)=2\left(v_{1}-v_{2}\right)$. Since $v_{i}+q c \in[1,3 / 2]$ and $v_{i}+(q+1) c \in[1,5 / 2]$ it follows that $l\left(v_{1}+q c\right)-l\left(v_{2}+q c\right)=-6\left(v_{1}-v_{2}\right)$ and $l\left(v_{1}+(q+1) c\right)-l\left(v_{2}+\right.$ $(q+1) c) \geqslant-3\left(v_{1}-v_{2}\right)$. Moreover, $l\left(v_{1}+j c\right) \geqslant l\left(v_{2}+j c\right)$ for $j \geqslant q+2$. Hence $M_{c}\left(v_{1}\right)-M_{c}\left(v_{2}\right) \geqslant-\left(v_{1}-v_{2}\right)$. That is, $M_{c}\left(v_{1}\right) \geqslant M_{c}\left(v_{2}\right)$, which means that $M_{c}(v)$ is decreasing on the interval $[1-q c, 1-(q-1) c]$.

Proof of (d). Fix $c \in(3 / 4,1)$. Consider $v_{1}$ and $v_{2}$ satisfying $(7 / 4-2 c) \vee 0<v_{1}<$ $v_{2}<1-c$. Then $7 / 4<v_{i}+2 c<1+c<2$ and $v_{i}+3 c>5 / 2$ for $i=1,2$. Hence $l\left(v_{1}\right)-$ $l\left(v_{2}\right)=2\left(v_{1}-v_{2}\right), l\left(v_{1}+c\right)-l\left(v_{2}+c\right)=2\left(v_{1}-v_{2}\right), l\left(v_{1}+2 c\right)-l\left(v_{2}+2 c\right)=-3\left(v_{1}-v_{2}\right)$, and $l\left(v_{1}+j c\right)-l\left(v_{2}+j c\right)=0$ for $j \geqslant 3$. Hence $M_{c}\left(v_{1}\right)-M_{c}\left(v_{2}\right)=v_{1}-v_{2}<0$, which gives (d).

Proof of (e). Fix $c \in[1,3 / 2]$. Clearly $M_{c}(v)$ is decreasing on $[1, \infty)$. Let $v_{1}, v_{2} \in$ $[0,1]$ with $v_{1} \leqslant v_{2}$. Then $v_{i}+c \in[1,5 / 2]$ and $v_{i}+j c \geqslant 2$ for $i=1,2$ and $j \geqslant 2$. Hence $l\left(v_{1}\right)-l\left(v_{2}\right)=2\left(v_{1}-v_{2}\right), l\left(v_{1}+c\right)-l\left(v_{2}+c\right) \geqslant-3\left(v_{1}-v_{2}\right)$ and $M_{c}\left(v_{1}\right)-M_{c}\left(v_{2}\right) \geqslant$ $-\left(v_{1}-v_{2}\right) \geqslant 0$. Hence $M_{c}(v)$ is decreasing on $[0, \infty)$. This yields (e) as in case 4 of (a).

Proof of (f). Consider $v_{1}$ and $v_{2}$ satisfying $(5 / 2-2 c) \vee 0<v_{1}<v_{2}<1$. Then $M_{c}\left(v_{1}\right)=l\left(v_{1}\right)<l\left(v_{2}\right)=M_{c}\left(v_{2}\right)$. Hence we have (f).

Proposition 3.5. Let $0<c^{\prime}<c$.
(i) $I(c) \varsubsetneqq I\left(c^{\prime}\right)$ if $c / c^{\prime} \in \mathbb{N}$,
(ii) $I(c) \nsubseteq I\left(c^{\prime}\right)$ and $I(c) \nsupseteq I\left(c^{\prime}\right)$ if $c / c^{\prime} \notin \mathbb{N}$.

Proof. The fact that $I(c) \subseteq I\left(c^{\prime}\right)$ for $c / c^{\prime} \in \mathbb{N}$ is already proved in Proposition 3.3 (ii). To show (ii) and the strict inclusion in (i) we may and do assume $d=1$. (Indeed, by considering Lévy measures on $\mathbb{R}^{d}$ in polar form $\lambda(\mathrm{d} \xi) \nu_{\xi}(\mathrm{d} u)$ where $\lambda$ is a point measure, the general case follows easily from the case $d=1$.) Let $\mu_{0}$ be a Poisson distribution. Let $\mu_{1}=\Phi_{c^{\prime}}\left(\mu_{0}\right)$ and $\mu_{2}=\Phi_{c}\left(\mu_{0}\right)$. Then $Q\left(\mu_{1}\right)=\{0\} \cup\left\{c^{\prime} / n: n \in \mathbb{N}\right\}$
and $Q\left(\mu_{2}\right)=\{0\} \cup\{c / n: n \in \mathbb{N}\}$ by Proposition 3.4. Thus, $\mu_{1} \in I\left(c^{\prime}\right) \backslash I(c)$ and for $c / c^{\prime} \notin \mathbb{N}$ we have $\mu_{2} \in I(c) \backslash I\left(c^{\prime}\right)$.

Remark 3.6. Let $c_{n}$ and $c$ be nonnegative with $c_{n} \rightarrow c$. Assume $c_{n} \neq c$ for all $n$. Then, as a complement to Theorem 3.1, we have

$$
\begin{equation*}
S D\left(\mathbb{R}^{d}\right) \varsubsetneqq \bigcap_{n \geqslant 1} I\left(c_{n}\right) \varsubsetneqq I(c) . \tag{3.7}
\end{equation*}
$$

Here we let $I(0)=I D\left(\mathbb{R}^{d}\right)$ as in Remark 2.9. The first inclusion in (3.7) follows since $S D\left(\mathbb{R}^{d}\right) \subseteq I\left(c^{\prime}\right)$ for all $c^{\prime}$. To see that it is strict choose $\mu \in I D\left(\mathbb{R}^{d}\right)$ with $Q(\mu)=[0, c+k]$ where $k>0$ is such that $c_{n} \leqslant c+k$ for all $n$. Such a $\mu$ exists by a straightforward extension of Proposition 3.4 (iv). Hence, $\mu \in I\left(c_{n}\right)$ for all $n$ and $\mu \notin S D\left(\mathbb{R}^{d}\right)$.

Let $\mu \in \bigcap_{n \geqslant 1} I\left(c_{n}\right)$. Then $c_{n} \in Q(\mu)$ for all $n$. Since $Q(\mu)$ is closed by Proposition 3.3 it follows that $c \in Q(\mu)$. To see that the last inclusion is strict let us consider the cases $c=0$ and $c>0$ separately. If $c=0$ then $I(c)=I D\left(\mathbb{R}^{d}\right)$. Hence, choosing $\mu \in I D\left(\mathbb{R}^{d}\right)$ such that $\mu \notin I\left(c^{\prime}\right)$ for any $c^{\prime}>0$ it follows that the inclusion is strict. In the case $c>0$ we can take $\mu \in I D\left(\mathbb{R}^{d}\right)$ such that $Q(\mu)=\{0\} \cup\{c / n: n=1,2, \ldots\}$. Then $\mu \in I(c) \backslash \bigcap_{n \geqslant 1} I\left(c_{n}\right)$.

## 4. Relations to other classes of distributions

We use the following notation as in [10]. For $\alpha \in(0,2]$ let $\mathfrak{S}_{\alpha}\left(\mathbb{R}^{d}\right)$ be the set of $\alpha$-stable distributions on $\mathbb{R}^{d}$. That is, a distribution $\mu$ on $\mathbb{R}^{d}$ is in $\mathfrak{S}_{\alpha}\left(\mathbb{R}^{d}\right)$ if and only if for every $n \in \mathbb{N}$ there is a $\beta \in \mathbb{R}^{d}$ such that $\widehat{\mu}(z)^{n}=\widehat{\mu}\left(n^{1 / \alpha} z\right) e^{i\langle\beta, z\rangle}$. (The definition of stability of $\mu$ is different from [12], p. 76, when $\mu$ is trivial.) Then $\mathfrak{S}\left(\mathbb{R}^{d}\right)=\bigcup_{\alpha \in(0,2]} \mathfrak{S}_{\alpha}\left(\mathbb{R}^{d}\right)$ is the set of stable distributions. For any distinct $\alpha$ and $\alpha^{\prime}$ in $(0,2] \mathfrak{S}_{\alpha}\left(\mathbb{R}^{d}\right) \cap \mathfrak{S}_{\alpha^{\prime}}\left(\mathbb{R}^{d}\right)$ is exactly the class of trivial distributions. Let $\mathfrak{S}_{\alpha}^{0}\left(\mathbb{R}^{d}\right)$ be the set of distributions $\mu$ on $\mathbb{R}^{d}$ such that for every $n \in \mathbb{N} \widehat{\mu}(z)^{n}=\widehat{\mu}\left(n^{1 / \alpha} z\right)$. Then $\mathfrak{S}^{0}\left(\mathbb{R}^{d}\right)=\bigcup_{\alpha \in(0,2]} \mathfrak{S}_{\alpha}^{0}\left(\mathbb{R}^{d}\right)$ is the set of strictly stable distributions. For any distinct $\alpha$ and $\alpha^{\prime}$ in $(0,2] \mathfrak{S}_{\alpha}^{0}\left(\mathbb{R}^{d}\right) \cap \mathfrak{S}_{\alpha^{\prime}}^{0}\left(\mathbb{R}^{d}\right)$ consists of the single element $\delta_{0}$. Distributions $\delta_{c}$ with $c \neq 0$ belong to $\mathfrak{S}_{1}^{0}\left(\mathbb{R}^{d}\right)$ but does not belong to any $\mathfrak{S}_{\alpha}^{0}\left(\mathbb{R}^{d}\right)$ with $\alpha \neq 1$.

Let $\left\{X_{s}\right\}_{s \in[0,1]}$ be a periodic Ornstein-Uhlenbeck process with parameter $c>0$ and background driving Lévy process $\left\{Z_{s}\right\}_{s \in[0,1]}$. Consider the following problem: To what extent are stability and the class $L_{m}$-property inherited from $\mathcal{L}\left(Z_{1}\right)$ to the stationary distribution $\mathcal{L}\left(X_{s}\right)$ ? We can recast this as the question to what extent $\Phi_{c}$ preserves stability and the class $L_{m}$-property.

Proposition 4.1. Let $c \in \mathbb{R} \backslash\{0\}, \mu_{0} \in I D\left(\mathbb{R}^{d}\right)$ and let $\mu=\Phi_{c}\left(\mu_{0}\right)$. Let $\alpha \in(0,2]$. Then
(i) $\mu \in \mathfrak{S}_{\alpha}^{0}\left(\mathbb{R}^{d}\right)$ if and only if $\mu_{0} \in \mathfrak{S}_{\alpha}^{0}\left(\mathbb{R}^{d}\right)$;
(ii) $\mu \in \mathfrak{S}_{\alpha}\left(\mathbb{R}^{d}\right)$ if and only if $\mu_{0} \in \mathfrak{S}_{\alpha}\left(\mathbb{R}^{d}\right)$.

Proof. Let $\Psi_{0}$ be the characteristic exponent of $\widehat{\mu}_{0}$. (i) Assume $\mu_{0} \in \mathbb{S}_{\alpha}^{0}\left(\mathbb{R}^{d}\right)$. Then $a \Psi_{0}(z)=\Psi_{0}\left(a^{1 / \alpha} z\right)$ for all $a>0$. This implies

$$
\widehat{\mu}(z)^{a}=\exp \int_{0}^{1} a \Psi_{0}\left(e^{c s} z\right) \mathrm{d} s=\exp \int_{0}^{1} \Psi_{0}\left(e^{c s} a^{1 / \alpha} z\right) \mathrm{d} s=\widehat{\mu}\left(a^{1 / \alpha} z\right)
$$

Hence $\mu \in \mathfrak{S}_{\alpha}^{0}\left(\mathbb{R}^{d}\right)$. Conversely, if $\mu \in \mathfrak{S}_{\alpha}^{0}\left(\mathbb{R}^{d}\right)$ then $\widehat{\mu}(z)^{a}=\widehat{\mu}\left(a^{1 / \alpha} z\right)$ for all $a>0$. Then by Corollary 2.8 we have $a \Psi_{0}(z)=\Psi_{0}\left(a^{1 / \alpha} z\right)$, which implies $\mu_{0} \in \mathfrak{S}_{\alpha}^{0}\left(\mathbb{R}^{d}\right)$.

Assume $\mu_{0} \in \mathfrak{S}_{\alpha}\left(\mathbb{R}^{d}\right)$. Then for any $a>0$ there is $\gamma_{a} \in \mathbb{R}^{d}$ such that $a \Psi_{0}(z)=$ $\Psi_{0}\left(a^{1 / \alpha} z\right)+\mathrm{i}\left\langle\gamma_{a}, z\right\rangle$. Then

$$
\begin{aligned}
\widehat{\mu}(z)^{a} & =\exp \int_{0}^{1} a \Psi_{0}\left(e^{c s} z\right) \mathrm{d} s=\exp \left(\int_{0}^{1} \Psi_{0}\left(e^{c s} a^{1 / \alpha} z\right) \mathrm{d} s+\mathrm{i} \int_{0}^{1}\left\langle\gamma_{a}, e^{c s} z\right\rangle \mathrm{d} s\right) \\
& =\widehat{\mu}\left(a^{1 / \alpha} z\right) e^{\mathrm{i} c^{-1}\left(e^{c}-1\right)\left\langle\gamma_{a}, z\right\rangle}
\end{aligned}
$$

Hence $\mu \in \mathfrak{S}_{\alpha}\left(\mathbb{R}^{d}\right)$ and conversely.
Let us recall the definition of the classes $L_{m}\left(\mathbb{R}^{d}\right)$ for $m=0,1, \ldots$. Let $L_{0}\left(\mathbb{R}^{d}\right)=$ $S D\left(\mathbb{R}^{d}\right)$. That is, $\mu$ is in $L_{0}\left(\mathbb{R}^{d}\right)$ if for every $b>1$ there is a distribution $\mu_{b}$ such that

$$
\begin{equation*}
\widehat{\mu}(z)=\widehat{\mu}\left(b^{-1} z\right) \widehat{\mu}_{b}(z) . \tag{4.1}
\end{equation*}
$$

Sometimes $\mu_{b}$ is referred to as an innovation distribution generated by $\mu$. For $m=$ $1,2, \ldots, L_{m}\left(\mathbb{R}^{d}\right)$ is the class of distributions $\mu$ such that, for every $b>1$, there exists $\mu_{b} \in L_{m-1}\left(\mathbb{R}^{d}\right)$ satisfying (4.1). We have

$$
\begin{equation*}
L_{0}\left(\mathbb{R}^{d}\right) \supsetneqq L_{1}\left(\mathbb{R}^{d}\right) \supsetneqq L_{2}\left(\mathbb{R}^{d}\right) \supsetneqq \cdots \supsetneqq L_{\infty}\left(\mathbb{R}^{d}\right) \supsetneqq \mathfrak{S}\left(\mathbb{R}^{d}\right), \tag{4.2}
\end{equation*}
$$

where $L_{\infty}\left(\mathbb{R}^{d}\right)=\bigcap_{0 \leq m<\infty} L_{m}\left(\mathbb{R}^{d}\right)$. These sets were introduced by Urbanik [15], [16] and further developed by Sato [11] in connection to limit distribution theory for sums of independent random vectors.

Remark 4.2. We recall a convenient characterization of $L_{m}\left(\mathbb{R}^{d}\right)$.
(i) Let $f: \mathbb{R} \longrightarrow \mathbb{R}$. For $\epsilon>0$ let $\Delta_{\epsilon}$ be the difference operator $\Delta_{\epsilon} f(v)=f(v+\epsilon)-$ $f(v)$, and $\Delta_{\epsilon}^{n}$ be the $n$th iteration of $\Delta_{\epsilon}$. Hence,

$$
\Delta_{\epsilon}^{n} f(v)=\sum_{j=0}^{n}(-1)^{n-j}\binom{n}{j} f(v+j \epsilon) .
$$

Let $\Delta_{\epsilon}^{0} f=f$. We say that $f$ is monotone of order $n$ if $\Delta_{\epsilon}^{j} f \geqslant 0$ for all $\epsilon>0$ and $j=0,1, \ldots, n$. From [10], Lemma 18, or [11], Lemma 3.2, we know that $f$ is monotone of order $n \geqslant 2$ if and only if $f \in C^{n-2}, f^{(j)} \geqslant 0$ for $j=0,1, \ldots, n-2$, and $f^{(n-2)}$ is increasing and convex.
(ii) Let $\mu \in I D\left(\mathbb{R}^{d}\right)$ have Lévy measure $\nu$ with polar decomposition $\lambda(\mathrm{d} \xi) \nu_{\xi}(\mathrm{d} u)$. By [12], Theorem 15.8, $\mu \in L_{0}\left(\mathbb{R}^{d}\right)=S D\left(\mathbb{R}^{d}\right)$ if and only if $\nu\left(b^{-1} B\right) \geqslant \nu(B)$ for $B \in \mathcal{B}\left(\mathbb{R}^{d}\right)$ and $b>1$. This is the case if and only if, for $\lambda$-a.e. $\xi, \nu_{\xi}\left(b^{-1} B\right) \geqslant \nu_{\xi}(B)$ for $B \in \mathcal{B}((0, \infty))$ and $b>1$. The latter is equivalent to $\nu_{\xi}\left(b^{-1} a_{1}, b^{-1} a_{2}\right] \geqslant \nu_{\xi}\left(a_{1}, a_{2}\right]$ for $0<a_{1}<a_{2}$ and $b>1$, which is equivalent to the property that $\nu_{\xi}\left(e^{-v-\epsilon}, e^{-v}\right]$ is increasing in $v$ for every $\epsilon>0$. If $\mu \in L_{0}\left(\mathbb{R}^{d}\right)$ then $\nu_{\xi}(\mathrm{d} u)=k_{\xi}(u) \mathrm{d} u / u$ with $k_{\xi}$ being decreasing, and conversely.
(iii) Let $\mu \in L_{0}\left(\mathbb{R}^{d}\right)$. That is, $\mu$ is selfdecomposable. Hence, the Lévy measure of $\mu$ is decomposed in polar form as $\lambda(\mathrm{d} \xi) k_{\xi}(u) \mathrm{d} u / u$, where $k_{\xi}(u)$ is decreasing. Let $h_{\xi}(v)=k_{\xi}\left(e^{-v}\right)$. Let $m \in\{1,2, \ldots\}$. From [10], Theorem 20, or [11], Theorems 3.2 and 3.3, we know that $\mu$ is in $L_{m}\left(\mathbb{R}^{d}\right)$ if and only if $h_{\xi}$ is monotone of order $m+1$ for $\lambda$-a.e. $\xi$.

Remark 4.3. Let $c>0$. For $\mu_{0} \in I D_{\log }\left(\mathbb{R}^{d}\right)$ we define $\widetilde{\Phi}_{-c}\left(\mu_{0}\right)=\mathcal{L}\left(\int_{0}^{\infty} e^{-c t} \mathrm{~d} Z_{t}\right)$ where $\left\{Z_{t}\right\}_{t \geqslant 0}$ is a Lévy process with $\mathcal{L}\left(Z_{1}\right)=\mu_{0}$. Then $\widetilde{\Phi}_{-c}$ is a one-to-one mapping from $I D_{\log }\left(\mathbb{R}^{d}\right)$ onto $L_{0}\left(\mathbb{R}^{d}\right)$, as is explained in Section 1 in connection to usual Ornstein-Uhlenbeck processes driven by Lévy processes. Let $m \in\{-1,0,1, \ldots\}$ and $j \in\{1,2, \ldots\}$. The following facts are known: the class of $\mu_{0}$ for which the $j$ th iteration $\widetilde{\Phi}_{-c}^{j}$ of $\widetilde{\Phi}_{-c}$ is definable coincides with $I D_{\log ^{j}}\left(\mathbb{R}^{d}\right) ; \mu \in L_{m+j}\left(\mathbb{R}^{d}\right)$ if and only if $\mu=\widetilde{\Phi}_{-c}^{j}\left(\mu_{0}\right)$ with some $\mu_{0} \in L_{m}\left(\mathbb{R}^{d}\right) \cap I D_{\log ^{j}}\left(\mathbb{R}^{d}\right)$ where we let $L_{-1}\left(\mathbb{R}^{d}\right)=I D\left(\mathbb{R}^{d}\right)$. See [5], [10] (Theorems 46 and 49), [14].

Proposition 4.4. Let $M \subseteq I D\left(\mathbb{R}^{d}\right)$ and $c \in \mathbb{R} \backslash\{0\}$. Assume that $M$ is completely closed and satisfies $(\mathrm{P})$. Then $\Phi_{c}(M)=\Phi_{-c}(M) \subseteq M$ and $\Phi_{c}(M)$ is completely closed and satisfies $(\mathrm{P})$. Moreover, we have $\Phi_{n c}(M) \subseteq \Phi_{c}(M)$ for all $n \in \mathbb{N}$. In addition, if $M$ satisfies ( D$)$ then $\Phi_{c}(M)$ satisfies $(\mathrm{D})$; if $M$ satisfies (GJ) then $\Phi_{c}(M)$ satisfies (GJ).

Proof. Assume $M$ is completely closed and satisfies (P). Let $\mu_{0} \in M$. Then $\Phi_{c}\left(\mu_{0}\right)=$ $\mathcal{L}\left(\int_{0}^{1} e^{c s} \mathrm{~d} Z_{s}\right)$, where $\left\{Z_{s}\right\}_{s \geqslant 0}$ is a Lévy process with $\mathcal{L}\left(Z_{1}\right)=\mu_{0}$. There is a sequence
of positive step functions $f_{1}, f_{2}, \ldots$ such that

$$
\int_{0}^{1} e^{c s} \mathrm{~d} Z_{s}=\lim _{n \rightarrow \infty} \int_{0}^{1} f_{n}(s) \mathrm{d} Z_{s}
$$

where the limit is in probability. There are $0=s_{n, 0}<s_{n, 1}<\ldots<s_{n, k}=1$ and $a_{n, j}>0$ such that

$$
\int_{0}^{1} f_{n}(s) \mathrm{d} Z_{s}=\sum_{j=1}^{k} a_{n, j}\left(Z_{s_{n, j}}-Z_{s_{n, j-1}}\right) .
$$

The distribution of the right-hand side belongs to $M$ by (C1), (C3) and (P). Hence $\Phi_{c}\left(\mu_{0}\right) \in M$ by $(\mathrm{C} 2)$. This proves $\Phi_{c}(M) \subseteq M$.

To show $\Phi_{c}(M)=\Phi_{-c}(M)$, let $\mu=\Phi_{c}\left(\mu_{0}\right)$ with $\mu_{0} \in M$. Let $\mu^{\prime}=\Phi_{-c}\left(\mu_{0}\right)$. By Proposition 2.4 (iv), $\widehat{\mu^{\prime}}(z)=\widehat{\mu}\left(e^{-c} z\right)$, that is, $\widehat{\mu^{\prime}}\left(e^{c} z\right)=\widehat{\mu}(z)$. Hence $\mu \in \Phi_{-c}(M)$ by $(\mathrm{C} 3)$ and $\Phi_{c}(M) \subseteq \Phi_{-c}(M)$. Changing the role of $c$ and $-c$, we get the converse inclusion.

It follows from Proposition 2.10 that $\Phi_{c}(M)$ satisfies (C2). It is moreover easily verified that $\Phi_{c}(M)$ satisfies (C1) and (C3), so let us verify ( P ) for $\Phi_{c}(M)$. Let $\mu=\mathcal{L}\left(\int_{0}^{1} e^{c s} \mathrm{~d} Z_{s}\right)$, where $\mathcal{L}\left(Z_{1}\right) \in M$. We have $\mu^{t}=\mathcal{L}\left(\int_{0}^{1} e^{c s} \mathrm{~d}_{s} Z_{t s}\right)$. Since $\mathcal{L}\left(Z_{t s}\right) \in M$ for all $t, s \geqslant 0$ by (P) of $M$ it follows that $\mu^{t} \in \Phi_{c}(M)$.

Assume $\mu \in \Phi_{n c}(M)$ and let $\left\{Z_{s}\right\}_{s \geqslant 0}$ be a Lévy process with $\mu=\mathcal{L}\left(\int_{0}^{1} e^{n c s} \mathrm{~d} Z_{s}\right)$. By Remark 2.2 we have

$$
\begin{aligned}
& \int_{0}^{1} e^{n c s} \mathrm{~d} Z_{s}=\int_{0}^{n} e^{c s} \mathrm{~d}_{s} Z_{s / n}=\sum_{j=0}^{n-1} \int_{j}^{j+1} e^{c s} \mathrm{~d}_{s} Z_{s / n} \\
= & \sum_{j=0}^{n-1} e^{j c} \int_{j}^{j+1} e^{c(s-j)} \mathrm{d}_{s} Z_{s / n} \stackrel{\mathrm{~d}}{=} \sum_{j=0}^{n-1} e^{j c} \int_{0}^{1} e^{c s} \mathrm{~d}_{s} Z_{s / n}^{j},
\end{aligned}
$$

where $\left\{Z_{s}^{j}\right\}_{s \geqslant 0}, j=0, \ldots, n-1$, are independent copies of $\left\{Z_{s}\right\}_{s \geqslant 0}$ and $\stackrel{\text { d }}{=}$ denotes equality in distribution. We have $\mathcal{L}\left(Z_{u}\right) \in M$ for all $u \geqslant 0$ by (P) of $M$. Hence $\mathcal{L}\left(\int_{0}^{1} e^{c s} \mathrm{~d} Z_{s / n}^{j}\right) \in \Phi_{c}(M)$ and thus $\mu=\mathcal{L}\left(\sum_{j=0}^{n-1} e^{j c} \int_{0}^{1} e^{c s} \mathrm{~d} Z_{s / n}^{j}\right) \in \Phi_{c}(M)$ since $\Phi_{c}(M)$ is completely closed.

If $M$ satisfies (D) then it is easy to see that $\Phi_{c}(M)$ satisfies (D).
Assume that $M$ satisfies (GJ). Let $\mu=\Phi_{c}\left(\mu_{0}\right)$ with $\mu_{0} \in M$ and $\mu=\mu_{G} * \mu_{J}$ where $\mu_{G}$ and $\mu_{J}$ are, respectively, Gaussian and purely non-Gaussian. We have $\mu_{0}=\mu_{0 G} * \mu_{0 J}$ with $\mu_{0 G}$ Gaussian and $\mu_{0 J}$ purely non-Gaussian. Then $\mu=\Phi_{c}\left(\mu_{0}\right)=$ $\Phi_{c}\left(\mu_{0 G}\right) * \Phi_{c}\left(\mu_{0 J}\right)$ by Proposition 2.4 (ii). By Proposition $2.3 \Phi_{c}\left(\mu_{0 G}\right)$ is Gaussian and $\Phi_{c}\left(\mu_{0 J}\right)$ is purely non-Gaussian. It follows that there is $b \in \mathbb{R}^{d}$ such that $\mu_{G}=$
$\Phi_{c}\left(\mu_{0 G}\right) * \delta_{b}$ and $\mu_{J}=\Phi_{c}\left(\mu_{0 J}\right) * \delta_{-b}$. Since $\mu_{0 G}, \mu_{0 J} \in M$ by (GJ) of $M$ we get $\mu_{G}, \mu_{J} \in \Phi_{c}(M)$ by (C3) of $\Phi_{c}(M)$.

As an immediate consequence of Proposition 4.4 we have the following.
Corollary 4.5. Let $c \in \mathbb{R} \backslash\{0\}$, $m=0,1, \ldots, \infty$ and $n \in \mathbb{N}$. Then $\Phi_{c}\left(L_{m}\left(\mathbb{R}^{d}\right)\right)$ is completely closed and satisfies $(\mathrm{D}),(\mathrm{P})$ and $(\mathrm{GJ})$, and $\Phi_{n c}\left(L_{m}\left(\mathbb{R}^{d}\right)\right) \subseteq \Phi_{c}\left(L_{m}\left(\mathbb{R}^{d}\right)\right)=$ $\Phi_{-c}\left(L_{m}\left(\mathbb{R}^{d}\right)\right) \subseteq L_{m}\left(\mathbb{R}^{d}\right)$.

Indeed, $L_{m}\left(\mathbb{R}^{d}\right)$ is completely closed and satisfies (D), (P) and (GJ).
Theorem 4.6. Let $c \in \mathbb{R} \backslash\{0\}$. Then,
(i) $L_{0}\left(\mathbb{R}^{d}\right) \varsubsetneqq I(c)$;
(ii) $L_{m+1}\left(\mathbb{R}^{d}\right) \varsubsetneqq \Phi_{c}\left(L_{m}\left(\mathbb{R}^{d}\right)\right) \varsubsetneqq L_{m}\left(\mathbb{R}^{d}\right)$ for $m=0,1, \ldots$;
(iii) $L_{\infty}\left(\mathbb{R}^{d}\right)=\Phi_{c}\left(L_{\infty}\left(\mathbb{R}^{d}\right)\right)$.

Remark 4.7. Using Proposition 4.4 and (ii) above it is immediate that $\Phi_{c}\left(L_{m+1}\left(\mathbb{R}^{d}\right)\right) \varsubsetneqq$ $\Phi_{c}^{2}\left(L_{m}\left(\mathbb{R}^{d}\right)\right) \varsubsetneqq \Phi_{c}\left(L_{m}\left(\mathbb{R}^{d}\right)\right)$ for $c \neq 0$ and $m \geqslant 0$. We do not know the relation between $\Phi_{c}^{2}\left(L_{m}\left(\mathbb{R}^{d}\right)\right)$ and $L_{m+1}\left(\mathbb{R}^{d}\right)$.

Proof of Theorem 4.6. The assertion (i) is by Theorem 3.1.
(ii) The second inclusion is by Corollary 4.5. Let $\mu \in L_{m+1}\left(\mathbb{R}^{d}\right)$ with $m \geqslant 0$. There exists by (i) a distribution $\mu_{0} \in I D\left(\mathbb{R}^{d}\right)$ such that $\mu=\Phi_{c}\left(\mu_{0}\right)$. We show that $\mu_{0} \in L_{m}\left(\mathbb{R}^{d}\right)$, which yields the first inclusion. Decompose the Lévy measure of $\mu$ as $\lambda(\mathrm{d} \xi) k_{\xi}(u) \mathrm{d} u / u$ in polar form and let $h_{\xi}(v)=k_{\xi}\left(e^{-v}\right)$ for $v \in \mathbb{R}$. Let $\nu_{0}$ denote the Lévy measure of $\mu_{0}$.

Step 1. First we prove $\mu_{0} \in L_{0}\left(\mathbb{R}^{d}\right)$. (This concludes the proof when $m=0$.) Define $H_{\xi}(v):=\sum_{j=1}^{\infty} k_{\xi}\left(e^{j c-v}\right)=\sum_{j=1}^{\infty} h_{\xi}(v-j c)$ for $v \in \mathbb{R}$. Since by Theorem 2.7 we have $H_{\xi}(v)=\left(c a_{\xi}\right)^{-1} \nu_{0 \xi}\left(e^{-v}, \infty\right)$ it follows that

$$
\begin{aligned}
& \sum_{j=1}^{\infty} \Delta_{\epsilon} h_{\xi}(v-j c)=\Delta_{\epsilon} H_{\xi}(v) \\
= & \left(c a_{\xi}\right)^{-1}\left(\nu_{0 \xi}\left(e^{-v-\epsilon}, \infty\right)-\nu_{0 \xi}\left(e^{-v}, \infty\right)\right)=\left(c a_{\xi}\right)^{-1} \nu_{0 \xi}\left(e^{-v-\epsilon}, e^{-v}\right] .
\end{aligned}
$$

This implies that $\nu_{0 \xi}\left(e^{-v-\epsilon}, e^{-v}\right]$ is increasing in $v$ for $\epsilon>0$, since $h_{\xi}(v)$ is increasing and convex by Remark 4.2. Hence $\mu_{0} \in L_{0}\left(\mathbb{R}^{d}\right)$ by the same remark.

Step 2. Assume $m \geqslant 1$. Since $\mu_{0} \in L_{0}\left(\mathbb{R}^{d}\right)$ the polar decomposition of $\nu_{0}$ is $\lambda_{0}(\mathrm{~d} \xi) k_{0 \xi}(u) \mathrm{d} u / u$, where the function $h_{0 \xi}(v):=k_{0 \xi}\left(e^{-v}\right)$ is nonnegative and increasing. By Remark 4.2, in order to show $\mu_{0} \in L_{m}\left(\mathbb{R}^{d}\right)$ it suffices to show that $h_{0 \xi}(v)$ is
monotone of order $m+1$. Since $\mu \in L_{m+1}\left(\mathbb{R}^{d}\right)$ the same remark shows that $h_{\xi}(v)$ is monotone of order $m+2$ and hence that $h_{\xi}^{\prime}(w)=\frac{\mathrm{d}}{\mathrm{d} w} h_{\xi}(w)$ exists and is monotone of order $m+1$. Define $H_{\xi}(v)$ as above. Since $h_{\xi}(-\infty)=0$ and $h_{\xi}^{\prime}(w) \geqslant 0$ we have

$$
H_{\xi}(v)=\sum_{j=1}^{\infty} \int_{-\infty}^{v-j c} h_{\xi}^{\prime}(w) \mathrm{d} w=\sum_{j=1}^{\infty} \int_{-\infty}^{v} h_{\xi}^{\prime}(w-j c) \mathrm{d} w=\int_{-\infty}^{v}\left(\sum_{j=1}^{\infty} h_{\xi}^{\prime}(w-j c)\right) \mathrm{d} w .
$$

Hence $H_{\xi}(v)$ is of class $C^{1}$ with $\frac{\mathrm{d}}{\mathrm{d} v} H_{\xi}(v)=\sum_{j=1}^{\infty} h_{\xi}^{\prime}(v-j c)$. On the other hand,

$$
H_{\xi}(v)=\left(c a_{\xi}\right)^{-1} \nu_{0 \xi}\left(e^{-v}, \infty\right)=\left(c a_{\xi}\right)^{-1} \int_{e^{-v}}^{\infty} k_{0 \xi}(u) \frac{\mathrm{d} u}{u}
$$

and thus $\frac{\mathrm{d}}{\mathrm{d} v} H_{\xi}(v)=\left(c a_{\xi}\right)^{-1} k_{0 \xi}\left(e^{-v}\right)=\left(c a_{\xi}\right)^{-1} h_{0 \xi}(v)$. It follows that $\left(c a_{\xi}\right)^{-1} h_{0 \xi}(v)=$ $\sum_{j=1}^{\infty} h_{\xi}^{\prime}(v-j c)$ and

$$
\left(c a_{\xi}\right)^{-1}\left(\Delta_{\epsilon}^{i} h_{0 \xi}\right)(v)=\sum_{j=1}^{\infty}\left(\Delta_{\epsilon}^{i} h_{\xi}^{\prime}\right)(v-j c) \geqslant 0 \text { for } i=0,1, \ldots, m+1,
$$

completing the proof of the first inclusion. We will show in Examples 4.8-4.9 below that the two inclusions are strict.
(iii) We have from the first inclusion in (ii) and from the one-to-one property of $\Phi_{c}$ that

$$
L_{\infty}\left(\mathbb{R}^{d}\right)=\bigcap_{0 \leqslant m<\infty} L_{m+1}\left(\mathbb{R}^{d}\right) \subseteq \bigcap_{0 \leqslant m<\infty} \Phi_{c}\left(L_{m}\left(\mathbb{R}^{d}\right)\right)=\Phi_{c}\left(L_{\infty}\left(\mathbb{R}^{d}\right)\right)
$$

On the other hand, $\Phi_{c}\left(L_{\infty}\left(\mathbb{R}^{d}\right)\right) \subseteq L_{\infty}\left(\mathbb{R}^{d}\right)$ by Corollary 4.5.

Example 4.8. We show that the second inclusion in Theorem 4.6 (ii) is strict. We may and do assume $d=1$ and using Corollary 4.5 it suffices to consider the case $c>0$. Example (i) shows that the inclusion is strict when $m=0$, while example (ii) applies to the case $m \geqslant 1$.
(i) We construct $\mu \in L_{0}(\mathbb{R})$ such that $\mu_{0}=\Phi_{c}^{-1}(\mu) \notin L_{0}(\mathbb{R})$. Let $\mu \in I D(\mathbb{R})$ have Lévy measure given by $\nu(\mathrm{d} u)=k(u) 1_{(0, \infty)}(u) \mathrm{d} u / u$, where

$$
k(u)= \begin{cases}1 & \text { for } 0<u \leqslant 1 \\ u^{-\alpha} & \text { for } u>1\end{cases}
$$

with $\alpha>0$. Then $h(v)=k\left(e^{-v}\right)$ equals 1 for $v \geqslant 0$ and $e^{\alpha v}$ for $v \leqslant 0$. This $h$ is increasing but not convex. Hence $\mu$ is in $L_{0}(\mathbb{R})$ but not in $L_{1}(\mathbb{R})$. Let $\mu_{0}=\Phi_{c}^{-1}(\mu)$, and denote the Lévy measure of $\mu_{0}$ by $\nu_{0}$. As in the proof of Theorem 4.6 we have $H(v+\epsilon)-H(v)=\frac{1}{c} \nu_{0}\left(e^{-v-\epsilon}, e^{-v}\right]$, where $H(v)=\sum_{j=1}^{\infty} h(v-j c)$. Since $H(v)=$ $e^{\alpha(v-c)} /\left(1-e^{-\alpha c}\right)$ for $v \leqslant c$ and $H(v)=1+e^{\alpha(v-2 c)} /\left(1-e^{-\alpha c}\right)$ for $c \leqslant v \leqslant 2 c$, we
can check that $H^{\prime}(c-)>H^{\prime}(c+)$, which shows that $H(v)$ is not convex. Thus for some $\epsilon>0 \nu_{0}\left(e^{-v-\epsilon}, e^{-v}\right]=c(H(v+\epsilon)-H(v))$ is not increasing in $v$. Therefore $\mu_{0} \notin L_{0}(\mathbb{R})$.
(ii) For $\mu \in L_{1}(\mathbb{R})$ we have $\mu_{0}=\Phi_{c}^{-1}(\mu) \in L_{0}(\mathbb{R})$ by Theorem 4.6. Let the Lévy measures of $\mu$ and $\mu_{0}$ be, respectively, $\nu(\mathrm{d} u)=1_{(0, \infty)}(u) k(u) \mathrm{d} u / u$ and $\nu_{0}(\mathrm{~d} u)=$ $1_{(0, \infty)}(u) k_{0}(u) \mathrm{d} u / u$. Defining $h(v)=k\left(e^{-v}\right)$ and $h_{0}(v)=k_{0}\left(e^{-v}\right)$, Theorem 2.7 shows that $h(v)=c^{-1} \int_{v}^{v+c} h_{0}(w) \mathrm{d} w$. Hence

$$
\begin{equation*}
h^{\prime}(v)=c^{-1}\left(h_{0}(v+c)-h_{0}(v)\right) \tag{4.3}
\end{equation*}
$$

if $h_{0}$ is continuous.
Let $m \geqslant 1$. We construct $\mu_{0} \in L_{0}(\mathbb{R})$ such that
(a) $h_{0}(v)$ is continuous and monotone of order $m$ but not of order $m+1$,
(b) $h_{0}(v+c)-h_{0}(v)$ is monotone of order $m$,
which by (4.3) and Remark 4.2 shows that $\mu_{0} \in L_{m-1}(\mathbb{R}) \backslash L_{m}(\mathbb{R})$ and $\mu=\Phi_{c}\left(\mu_{0}\right) \in$ $L_{m}(\mathbb{R})$.

Let $0<a<2, a<b$ and $b \leqslant a e^{a c}$. Define $g_{j}(v)$ for $j=0, \ldots, m-1$ as

$$
\begin{aligned}
& g_{0}(v)= \begin{cases}e^{b v} & \text { for } v \leqslant 0 \\
e^{a v} & \text { for } v \geqslant 0,\end{cases} \\
& g_{j}(v)=\int_{-\infty}^{v} g_{j-1}(w) \mathrm{d} w, \quad \text { for } j=1, \ldots, m-1 .
\end{aligned}
$$

Using that $0<a<2$ and $b>0$, we see that $h_{0}(v):=g_{m-1}(v)$ satisfies the integrability condition

$$
\int_{-\infty}^{0} h_{0}(v) \mathrm{d} v+\int_{0}^{\infty} e^{-2 v} h_{0}(v) \mathrm{d} v<\infty
$$

in [11], p. 215, which implies that $\nu_{0}$ constructed from $h_{0}(v)$ is indeed a Lévy measure. Since $g_{0}(v)$ is nonnegative, continuous and increasing but not convex, we have (a) by Remark 4.2. Using that $0<b \leqslant a e^{a c}$ and $a<b$, it is readily verified that $g_{0}(v+c)-g_{0}(v)$ is increasing and nonnegative. When $m=1$, this fact gives (b). Assume $m \geqslant 2$. Let $f(v)=h_{0}(v+c)-h_{0}(v)$ and notice that

$$
\begin{equation*}
f^{(j)}(v)=\int_{-\infty}^{v}\left(g_{m-2-j}(w+c)-g_{m-2-j}(w)\right) \mathrm{d} w, \quad j=0, \ldots, m-2 \tag{4.4}
\end{equation*}
$$

By (4.4) this implies that $f^{(m-2)}$ is increasing, nonnegative and convex, which in turn implies that $f^{(j)}, j=0, \ldots, m-2$, have the same properties. Now (b) follows from Remark 4.2.

Example 4.9. We show that the first inclusion in Theorem 4.6 (ii) is strict. It suffices to give an example of a distribution in $\Phi_{-c}\left(L_{m}\left(\mathbb{R}^{d}\right)\right) \backslash L_{m+1}\left(\mathbb{R}^{d}\right)$ in the case $c>0$ and $d=1$.

Let $b=e^{c}$. By Proposition 5.1 in Appendix there is $\mu \in L_{m+1}(\mathbb{R})$ such that $\widehat{\mu}(z)=\widehat{\mu}\left(b^{-1} z\right) \widehat{\mu}_{b}(z)$ with $\mu_{b} \in L_{m}(\mathbb{R}) \backslash L_{m+1}(\mathbb{R})$. By Remark 4.3, there is $\mu_{0} \in$ $L_{m}(\mathbb{R}) \cap I D_{\log }(\mathbb{R})$ such that $\mu=\widetilde{\Phi}_{-c}\left(\mu_{0}\right)$. Let $\left\{Z_{s}\right\}_{s \geqslant 0}$ be a Lévy process with $\mathcal{L}\left(Z_{1}\right)=\mu_{0}$. Since

$$
\int_{0}^{\infty} e^{-c s} \mathrm{~d} Z_{s}=\int_{1}^{\infty} e^{-c s} \mathrm{~d} Z_{s}+\int_{0}^{1} e^{-c s} \mathrm{~d} Z_{s} \text { and } \int_{1}^{\infty} e^{-c s} \mathrm{~d} Z_{s} \stackrel{\mathrm{~d}}{=} e^{-c} \int_{0}^{\infty} e^{-c s} \mathrm{~d} Z_{s}
$$

we have $\mu_{b}=\Phi_{-c}\left(\mu_{0}\right)$. Hence $\Phi_{-c}\left(L_{m}(\mathbb{R})\right)$ contains $\mu_{b} \in L_{m}(\mathbb{R}) \backslash L_{m+1}(\mathbb{R})$.

Remark 4.10. For $m \in\{0,1, \ldots, \infty\}$ let $M_{m}$ be the class of probability measures $\rho$ on $\mathbb{R}^{d}$ such that

$$
\begin{equation*}
\widehat{\rho}(z)=\frac{\widehat{\mu}(z)}{\widehat{\mu}\left(b^{-1} z\right)} \tag{4.5}
\end{equation*}
$$

for some $\mu \in L_{m}\left(\mathbb{R}^{d}\right)$ and some $b>1$. Thus, $\rho$ is an innovation distribution generated by a selfdecomposable distribution $\mu \in L_{m}\left(\mathbb{R}^{d}\right)$. Characterization of this class $M_{m}$ is of some interest. Let us see that

$$
M_{m}=\bigcup_{c>0} \Phi_{c}\left(L_{m-1}\left(\mathbb{R}^{d}\right)\right) \cap I D_{\log }\left(\mathbb{R}^{d}\right)
$$

where $L_{-1}\left(\mathbb{R}^{d}\right)=I D\left(\mathbb{R}^{d}\right)$ and $L_{\infty-1}\left(\mathbb{R}^{d}\right)=L_{\infty}\left(\mathbb{R}^{d}\right)$. Indeed, let $\rho \in M_{m}$ be given by (4.5) with $\mu \in L_{m}\left(\mathbb{R}^{d}\right)$. Let $c=\log b$. Then $\mu=\widetilde{\Phi}_{-c}\left(\mu_{0}\right)$ for some $\mu_{0} \in I D_{\log }\left(\mathbb{R}^{d}\right) \cap$ $L_{m-1}\left(\mathbb{R}^{d}\right)$ by Remark 4.3. Repeating the arguments in Example 4.9 it follows that $\rho=\Phi_{-c}\left(\mu_{0}\right) \in \Phi_{-c}\left(L_{m-1}\left(\mathbb{R}^{d}\right)\right)=\Phi_{c}\left(L_{m-1}\left(\mathbb{R}^{d}\right)\right)$. Since $\mu_{0} \in I D_{\log }\left(\mathbb{R}^{d}\right)$ we have also $\rho \in I D_{\log }\left(\mathbb{R}^{d}\right)$ by Remark 2.12 (iv).

For the converse, let $\rho \in \Phi_{c}\left(L_{m-1}\left(\mathbb{R}^{d}\right)\right) \cap I D_{\log }\left(\mathbb{R}^{d}\right)$ with some $c>0$. Noting that $\Phi_{c}\left(L_{m-1}\left(\mathbb{R}^{d}\right)\right)=\Phi_{-c}\left(L_{m-1}\left(\mathbb{R}^{d}\right)\right)$, let $\mu_{0}=\Phi_{-c}^{-1}(\rho)$. Then $\mu_{0} \in L_{m-1}\left(\mathbb{R}^{d}\right)$ and since $\rho \in I D_{\log }\left(\mathbb{R}^{d}\right)$ we have also $\mu_{0} \in I D_{\log }\left(\mathbb{R}^{d}\right)$ by Remark 2.12 (iv). Let $\mu=\widetilde{\Phi}_{-c}\left(\mu_{0}\right)$. Then $\mu \in L_{m}\left(\mathbb{R}^{d}\right)$ by Remark 4.3 and $\widehat{\mu}(z)=\widehat{\rho}(z) \widehat{\mu}\left(e^{-c} z\right)$. Therefore $\rho \in M_{m}$.

## 5. Appendix

Recall that, for $m=0,1, \ldots, \mu \in L_{m}\left(\mathbb{R}^{d}\right)$ if and only if there exists for every $b>1$ a distribution $\mu_{b} \in L_{m-1}\left(\mathbb{R}^{d}\right)$ such that we have (4.1). (Here we let $L_{-1}\left(\mathbb{R}^{d}\right)=$ $\left.I D\left(\mathbb{R}^{d}\right)\right)$.

Proposition 5.1. Let $m \in\{0,1, \ldots\}$.
(a) There is a distribution $\mu \in L_{m}\left(\mathbb{R}^{d}\right) \backslash L_{m+1}\left(\mathbb{R}^{d}\right)$ such that, for every $b>1$, $\mu_{b} \in L_{m-1}\left(\mathbb{R}^{d}\right) \backslash L_{m}\left(\mathbb{R}^{d}\right)$.
(b) Let $b>1$. Then there exists $\mu \in L_{m}\left(\mathbb{R}^{d}\right) \backslash L_{m+1}\left(\mathbb{R}^{d}\right)$ such that $\mu_{b} \in L_{m}\left(\mathbb{R}^{d}\right)$.

Proof. Assume without loss of generality $d=1$. Let $\mu \in L_{0}(\mathbb{R})$ have Lévy measure $\nu$ given by $\nu(\mathrm{d} u)=1_{(0, \infty)}(u) k(u) \mathrm{d} u / u$. Let $h(v)=k\left(e^{-v}\right)$ for $v \in \mathbb{R}$. For $b>1$ the Lévy measure $\nu_{b}$ of $\mu_{b}$ is $\nu_{b}(\mathrm{~d} u)=1_{(0, \infty)}(u) k_{b}(u) \mathrm{d} u / u$, where $k_{b}(u)=k(u)-k(b u)$. Letting $h_{b}(v)=k_{b}\left(e^{-v}\right)$ we thus have $h_{b}(v)=h(v)-h(v-\log b)$.
(a) Let $\mu$ be the exponential distribution with mean 1. By [12], p. 45, we have $h(v)=e^{-e^{-v}}$. Hence $h^{\prime \prime}(v)=e^{-e^{-v}}\left(e^{-2 v}-e^{-v}\right)<0$ for $v>0$. Since

$$
h_{b}^{\prime}(v)=h^{\prime}(v)-h^{\prime}(v-\log b)=\int_{v-\log b}^{v} h^{\prime \prime}(r) \mathrm{d} r<0 \quad \text { if } v-\log b>0
$$

it follows from Remark 4.2 that $\mu_{b} \in I D(\mathbb{R}) \backslash L_{0}(\mathbb{R})$ for every $b>1$. Thus $\mu \in$ $L_{0}(\mathbb{R}) \backslash L_{1}(\mathbb{R})$. We have $\mu \in I D_{\log ^{m}}(\mathbb{R})$ and hence $\mu_{b} \in I D_{\log ^{m}}(\mathbb{R})$ for all $m$. Fix $c>0$. Let $\sigma=\widetilde{\Phi}_{-c}^{m}(\mu)$ and $\sigma_{b}=\widetilde{\Phi}_{-c}^{m}\left(\mu_{b}\right)$. Then $\widehat{\sigma}(z)=\widehat{\sigma}\left(b^{-1} z\right) \widehat{\sigma}_{b}(z)$. Observe that $\sigma \in L_{m}(\mathbb{R}) \backslash L_{m+1}(\mathbb{R})$ and $\sigma_{b} \in L_{m-1}(\mathbb{R}) \backslash L_{m}(\mathbb{R})$ by Remark 4.3. As an alternative example, we can use the distribution with Lévy measure $1_{(0,1)}(u) \mathrm{d} u / u$ instead of the exponential distribution.
(b) Given $b>1$, consider $\mu \in I D(\mathbb{R}) \backslash L_{0}(\mathbb{R})$ such that $\widehat{\mu}(z)=\widehat{\mu}\left(b^{-1} z\right) \widehat{\mu}_{b}(z)$ with $\mu_{b} \in I D(\mathbb{R})$. Such a $\mu$ exists due to the existence of a non-selfdecomposable distribution $\mu$ that is semi-selfdecomposable with span b, as in [12], Example 15.9. Further we can assume $\mu \in \bigcap_{1 \leqslant m<\infty} I D_{\log ^{m}(\mathbb{R})}$ by truncating its Lévy measure or by considering a semi-stable distribution. Fix $c>0$ and let $\sigma=\widetilde{\Phi}_{-c}^{m+1}(\mu)$. We have $\widehat{\sigma}(z)=\widehat{\sigma}\left(b^{-1} z\right) \widehat{\sigma}_{b}(z)$ with $\sigma_{b}=\widetilde{\Phi}_{-c}^{m+1}\left(\mu_{b}\right)$. Then, by Remark 4.3, $\sigma \in L_{m}(\mathbb{R}) \backslash L_{m+1}(\mathbb{R})$ and $\sigma_{b} \in L_{m}(\mathbb{R})$.

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