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# JEAN JACOD <br> Victor Pérez-Abreu <br> On martingales which are finite sums of independent random variables with time dependent coefficients 

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# On martingales which are finite sums of independent random variables with time dependent coefficients 

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## 1 Introduction

We consider the following problem: for a positive integer $n \geq 1$, let $U_{1}, \ldots, U_{n}$ be $n$ independent, integrable, centered, non-degenerate random variables. We are looking for conditions on a family of $n$ càdlàg functions $f_{1}, \ldots, f_{n}$ on $\mathbb{R}_{+}$with $f_{i}(0)=0$, under which the following process:

$$
\begin{equation*}
X_{t}=\sum_{i=1}^{n} f_{i}(t) U_{i} \tag{1}
\end{equation*}
$$

is a martingale, with respect to its own filtration $\left(\mathcal{F}_{t}\right)_{t \geq 0}$.
This (apparently) simple problem has a general solution given in Section 1. However, the answer is not quite satisfactory, since for example it does not allow to recognize whether there is a unique (up to the obvious multiplication by constants and time-changes) set ( $f_{i}$ ) meeting our condition.

To get more insight, we specialize in Section 3 to the case where $n=2$ and (for the most interesting results) with $U_{1}$ and $U_{2}$ having the same law. In this very particular situation we are able to give a complete description of all martingales of the form (1). This description emphasizes the particular role played by the stable distributions.

For the case $n \geq 3$, we have been unable to provide any interesting result of the same kind as for $n=2$.

## 2 A general result

Here is a general theorem solving (in principle) our problem.
Theorem 1. The process $X$ is a martingale if and only if it satisfies the following:
Condition [M]: There are an integer $p, 0 \leq p \leq n$, and deterministic times $0=$ $T_{0}<T_{1}<\ldots<T_{p}<T_{p+1}=\infty$, and $p$ linearly independent vectors $a_{j}=\left(a_{j}^{i}\right)_{1 \leq i \leq n}$ in $\mathbb{R}^{n}$ (when $p \geq 1$ ), such that, with $V_{0}=0$ and $V_{j}=\sum_{1 \leq i \leq n} a_{j}^{i} U_{i}$ for $j \geq 1$,
(M1) $\left(V_{j}\right)_{0 \leq j \leq p}$ is a discrete-time martingale;
(M2) $X_{t}=\sum_{1 \leq j \leq p} V_{j} 1_{\left[T_{,}, T_{j+1}\right)}(t)$.
Before proving this theorem, we state some remarks on the conditions. First, Condition (M2) implies that $f_{i}(t)=\sum_{1 \leq j \leq p} a_{j}^{i} 1_{\left[T_{j}, T_{j+1}\right)}(t)$, because of the following property:

$$
\begin{equation*}
\alpha_{i}, \beta_{i} \in \mathbb{R}, \quad \sum_{i=1}^{n} \alpha_{i} U_{i}=\sum_{i=1}^{n} \beta_{i} U_{i} \quad \text { a.s. } \quad \Rightarrow \quad \alpha_{i}=\beta_{i} \quad \forall i . \tag{2}
\end{equation*}
$$

Second, Condition (M1) is obviously difficult to verify, except when $p=0$ (it is void) and $p=1$ (it is obvious because $V_{1}$ is centered). Below we give an equivalent condition based on the characteristic functions $\varphi_{i}$ of $U_{i}$. We recall that each function $\varphi_{i}$ is $C^{1}$ with $\varphi_{i}^{\prime}(0)=0$. Then, when $p \geq 2$, (M1) is equivalent to the following:

Condition (M'1). For all $1 \leq l \leq p-1$ and all $v_{j}$ in $\mathbb{R}$,

$$
\begin{equation*}
\sum_{i=1}^{n}\left(a_{l+1}^{i}-a_{l}^{i}\right) \varphi_{i}^{\prime}\left(\sum_{j=1}^{l} a_{j}^{i} v_{j}\right) \prod_{k \neq i} \varphi_{k}\left(\sum_{j=1}^{l} a_{j}^{k} v_{j}\right)=0 \tag{3}
\end{equation*}
$$

We observe that (3) is the same as $E\left(\left(V_{l+1}-V_{l}\right) \exp i \sum_{j=1}^{l} v_{j} V_{j}\right)=0$. When the $\varphi_{i}$ 's do not vanish (so $\varphi_{i}=\exp \psi_{i}$ with $\psi_{i}$ of class $C^{1}$ and $\psi_{i}^{\prime}(0)=0$ ) this condition is also equivalent to:

Condition ( $\mathbf{M}$ "1). For all $1 \leq l \leq p-1$ and all $v_{j}$ in $\mathbb{R}$,

$$
\begin{equation*}
\sum_{i=1}^{n}\left(a_{l+1}^{i}-a_{l}^{i}\right) \psi_{i}^{\prime}\left(\sum_{j=1}^{l} a_{j}^{i} v_{j}\right)=0 \tag{4}
\end{equation*}
$$

Proof. The sufficient condition is obvious. For the necessary condition, we suppose that $X$ is a martingale and let $F(t)$ be the vector with components $\left(f_{i}(t)\right)_{1 \leq i \leq n}$. Denote by $E_{t}$ the linear space spanned by $(F(s): s \leq t)$, let $d_{t}=\operatorname{dim}\left(E_{t}\right), T_{-1}=-1$, $T_{j}=\inf \left(t: d_{t} \geq j\right)$ for $0 \leq j \leq n$, and $T_{n+1}=\infty$. Thus $T_{-1}<0=T_{0} \leq T_{1} \leq \ldots \leq$ $T_{p}<T_{p+1}=\infty$ for some $0 \leq p \leq n$, and $d_{0}=0$.

Let $0 \leq i \leq p$ with $T_{i}<T_{i+1}$ and consider $s, t$ such that $T_{i}<s<t<T_{i+1}$. Then $E_{t}=E_{s}$ is spanned by the linearly independent vectors $F\left(s_{1}\right), \ldots, F\left(s_{i}\right)$ with $s_{j} \leq s$ (if $i=0$, then $E_{t}=E_{s}=\{0\}$ ). Therefore, $X_{s}$ and $X_{t}$ are $\sigma\left(X_{s_{1}}, \ldots, X_{s_{i}}\right)$-measurable and thus $\mathcal{F}_{s}=\mathcal{F}_{t}=\sigma\left(X_{s_{1}}, \ldots, X_{s_{t}}\right)$ (which is the trivial $\sigma$-field when $i=0$ ). The martingale property $E\left(X_{t} \mid \mathcal{F}_{s}\right)=X_{s}$ yields $X_{t}=X_{s}$ a.s., and (2) gives $F(s)=F(t)$. It follows that $F($.$) is constant on \left(T_{i}, T_{i+1}\right)$ as well as on $\left[T_{i}, T_{i+1}\right)$ by right-continuity. Thus

$$
\begin{equation*}
T_{i}<T_{i+1} \quad \Rightarrow \quad d_{r}=i \quad \forall r \in\left[T_{i}, T_{i+1}\right) . \tag{5}
\end{equation*}
$$

In fact $0<T_{1}<\ldots<T_{p}$; otherwise we would be in one of the following two situations:
a) $0=T_{j}<T_{j+1}$ for some $1 \leq j \leq p$, and therefore $d_{T_{j}}=d_{0}=0$, which contradicts (5);
b) $T_{i-1}<T_{i}=T_{j}<T_{j+1}$ for $i, j$ with $1 \leq i<j \leq p$, in which case $d_{r}=i-1$ on $\left[T_{i-1}, T_{j}\right)$ by (3). This implies that $d_{T_{j}} \leq i$; being also impossible since $d_{T_{j}} \geq j$.

Since $0<T_{1}<\ldots<T_{p}$ holds, we trivially have (M2) with $a_{j}=F\left(T_{j}\right)$. Finally, (M2) and the martingale property of $X$ yield (M1).

## 3 The case $\mathbf{n}=2$

Let $\varphi_{i}$ be the characteristic function of $U_{i}$, and when $\varphi_{i}$ never vanishes we use the notation $\varphi_{i}=\exp \psi_{i}$ without further comment. In this section we always assume that $n=2$.

Theorem 2. The process $X$ is a martingale if and only if it has one of the following two (mutually exclusive) representations:
a) For some $\alpha, \beta \in \mathbb{R}, S_{1}, S_{2} \in(0, \infty]$

$$
\begin{equation*}
X_{t}=\alpha U_{1} 1_{\left[S_{1}, \infty\right)}(t)+\beta U_{2} 1_{\left[S_{2}, \infty\right)}(t) \tag{6}
\end{equation*}
$$

b) For some $0<T_{1}<T_{2}<\infty, \alpha, \alpha^{\prime}, \gamma, \gamma^{\prime} \in \mathbb{R}^{*}$ with $\gamma \neq \gamma^{\prime}$ and

$$
\begin{gather*}
\varphi_{1}^{\prime}(v) \varphi_{2}(\gamma v)+\gamma^{\prime} \varphi_{1}(v) \varphi_{2}^{\prime}(\gamma v)=0 \quad \forall v \in \mathbb{R},  \tag{7}\\
X_{t}=\alpha\left(U_{1}+\gamma U_{2}\right) 1_{\left[T_{1}, \infty\right)}(t)+\alpha^{\prime}\left(U_{1}+\gamma^{\prime} U_{2}\right) 1_{\left[T_{2}, \infty\right)}(t) . \tag{8}
\end{gather*}
$$

Remark. Since the coefficients in (8) do not vanish, the form (8) is indeed symmetric in $\left(U_{1}, U_{2}\right)$. When $\varphi_{1}$ and $\varphi_{2}$ do not vanish, (7) is equivalent to $\psi_{1}^{\prime}(v)+\gamma^{\prime} \psi_{2}^{\prime}(\gamma v)=0$, which is the same as $\psi_{1}(v)+\frac{\gamma^{\prime}}{\gamma} \psi_{2}(\gamma v)=0$, which in turn is equivalent to

$$
\begin{equation*}
\varphi_{1}(v)=\varphi_{2}(\gamma v)^{-\gamma^{\prime} / \gamma} \quad \forall v \in \mathbb{R} . \tag{9}
\end{equation*}
$$

Proof. Sufficient condition: That (a) gives a martingale is obvious. Condition (b) implies (M2) with $a_{1}^{1}=\alpha, a_{1}^{2}=\alpha \gamma, a_{2}^{1}=\alpha^{\prime}+a_{1}^{1}, a_{2}^{2}=\alpha^{\prime} \gamma^{\prime}+a_{1}^{2}$ and then (7) gives (M'1).

Necessary condition: We assume (M'1) and (M2). If $T_{1}=\infty$, then (a) holds with $\alpha=\beta=0$ and $S_{i}$ arbitrary. If $T_{1}<T_{2}=\infty$, then (a) holds with $\alpha=a_{1}^{1}, \beta=a_{1}^{2}$ and $S_{1}=S_{2}=T_{1}$.

Suppose now that $T_{1}<T_{2}<\infty$. We have $a_{1} \neq 0$, and since both (a) and (b) are symmetric in ( $U_{1}, U_{2}$ ), without lost of generality we assume that $a_{1}^{1} \neq 0$. Let $\alpha=a_{1}^{1}$ and $\gamma=a_{1}^{2} / \alpha$ and write $a_{2}^{i}=a_{1}^{i}+\beta^{i}$. Then the linear independence between $a_{1}$ and $a_{2}$ gives

$$
\begin{equation*}
\beta^{2} \neq \gamma \beta^{1} \tag{10}
\end{equation*}
$$

while (M'1) is

$$
\begin{equation*}
\beta^{1} \varphi_{1}^{\prime}(v) \varphi_{2}(\gamma v)+\beta^{2} \varphi_{1}(v) \varphi_{2}^{\prime}(\gamma v)=0 \quad \forall v \in \mathbb{R} . \tag{11}
\end{equation*}
$$

We assume first that $\gamma=0$. Recalling that $\varphi_{i}(0)=1, \varphi_{i}^{\prime}(0)=0$ and $\varphi_{i}^{\prime}$ is not identically 0 in any neighborhood of 0 (because $P\left(U_{i}=0\right)<1$ ), (11) yields $\beta^{1}=0$, that is, we have (a) with $S_{1}=T_{1}, S_{2}=T_{2}, \beta=\beta^{2}$.

Next, assume that $\gamma \neq 0$. Then there exists $\theta \in \mathbb{R}^{*}$ with $\varphi_{1}^{\prime}(\theta) \neq 0, \varphi_{1}(\theta) \neq 0$ and $\varphi_{2}(\gamma \theta) \neq 0$. Suppose for the time being that $\varphi_{2}^{\prime}(\gamma \theta)=0$. Then (11) yields $\beta^{1}=0$ and since there is another $\theta^{\prime} \in \mathbb{R}^{*}$ with $\varphi_{1}\left(\theta^{\prime}\right) \neq 0$ and $\varphi_{2}\left(\gamma \theta^{\prime}\right) \neq 0$, we also have $\beta^{2}=0$, which contradicts (10). Thus $\varphi_{2}^{\prime}(\gamma \theta) \neq 0$ and (10) and (11) yield $\beta^{1} \neq 0$ and $\beta^{2} \neq 0$. Hence we have (b) with $\gamma^{\prime}=\beta^{2} / \beta^{1}$ and $\alpha^{\prime}=\beta^{1}$ (note that $\gamma \neq \gamma^{\prime}$ follows from (10), and (7) is the same as (11)).

When $U_{1}$ and $U_{2}$ are arbitrary, it seems there is not much more to say. From now on we concentrate on the case where $U_{1}={ }^{d} U_{2}$, i.e. $\varphi_{1}=\varphi_{2}=\varphi$. In this situation, the existence of a martingale $X$ of the form (b) above depends on the existence of constants $\gamma, \gamma^{\prime} \in \mathbb{R}^{*}$ with $\gamma \neq \gamma^{\prime}$ and

$$
\begin{equation*}
\varphi^{\prime}(v) \varphi(\gamma v)+\gamma^{\prime} \varphi(v) \varphi^{\prime}(\gamma v)=0 \quad \forall v \in \mathbb{R} \tag{12}
\end{equation*}
$$

Let $D$ denote the set of all $\gamma \in \mathbb{R}^{*}$ for which (12) holds for some $\gamma^{\prime} \in \mathbb{R}^{*}$ with $\gamma^{\prime} \neq \gamma$. If $\gamma \in D$ there is a unique $\gamma^{\prime}=\delta(\gamma)$ satisfying (12), because we have seen before that for each $\gamma \neq 0$ there is $v \in \mathbb{R}$ with $\varphi(v) \neq 0$ and $\varphi^{\prime}(\gamma v) \neq 0$.

Theorem 3. a) If $U_{1}$ is symmetric about 0 , then one of the following three cases is satisfied:
(Cs-1) $\quad D=\{-1,1\}$.
(Cs-2) $D=\left\{r^{n},-r^{n}: n \in \mathbb{Z}\right\}$ for some $r>1$ and $\varphi$ never vanishes.
(Cs-3) $D=\mathbb{R}^{*}$. This is the case if and only if $U_{1}$ is stable with index $\rho \in(1,2]$, i.e. $\quad \varphi(u)=\epsilon^{-a|u|^{\rho}}$ for some $a>0$.
b) If $U_{1}$ is not symmetric about 0 , we are in one of the following five situations:
(Ca-1) $D=\{1\}$.
(Ca-2) $D=\{-1,1\}$. This is the case if and only if $\varphi=\rho e^{\eta}$, where $\rho$ and $\eta$ are real-valued, $\eta(0)=0$, and $\eta$ is constant on each open interval on which $\rho$ (or $\varphi$ ) does not vanish (necessarily $\varphi$ vanishes somewhere, and $\eta$ is not identically 0 , otherwise we would be in the symmetric case).
(Ca-3) $D=\left\{r^{n}: n \in \mathbb{Z}\right\}$ for some $r>1$ and $\varphi$ never vanishes.
(Ca-4) $D=\left\{r^{n},-r^{n+1 / 2}: n \in \mathbb{Z}\right\}$ for some $r>1$ and $\varphi$ never vanishes.
(Ca-5) $D=(0, \infty)$. This is the case if and only if $U_{1}$ is asymmetric strictly stable with index $\rho \in(1,2)$, i.e., $\varphi(u)=\epsilon^{-a|u|^{\rho}(1+i b s i g n(u))}$ for some $a>0, b \neq 0$, $|b| \leq \tan \left(\frac{\pi}{2(2-\rho)}\right)$.
c) There is a constant $\theta \in(1,2]$ such that $\delta(\gamma)=-\gamma /|\gamma|^{\theta} \quad$ (so $\delta(1)=-1$, and $\delta(-1)=1$ if $-1 \in D$ ), and $\theta=\rho$ in cases (Cs-3) and (Ca-5).

Therefore the martingales $X$ of the form (8) are indeed represented as

$$
\begin{equation*}
X_{t}=\alpha\left(U_{1}+\gamma U_{2}\right) 1_{\left[T_{1}, \infty\right)}(t)+\alpha^{\prime}\left(U_{1}-\gamma U_{2} /|\gamma|^{\theta}\right) 1_{\left[T_{2}, \infty\right)}(t), \tag{13}
\end{equation*}
$$

where $\alpha, \alpha^{\prime} \in \mathbb{R}^{*}, \quad 0<T_{1}<T_{2}<\infty$, and $\gamma \in D$.
Remark. There are of course examples of variables satisfying (Cs-1) or (Cs-3) in the symmetrical case, (Ca-1) in the asymmetrical case. We presume that (Cs-2) and (Ca-3) are not empty, and believe that (Ca-2) is empty (but we have been unable to prove these facts).

Before giving the proof of Theorem 3 we present some useful lemmas. First we note that $\gamma=1$ and $\gamma^{\prime}=-1$ always satisfy (12), so $1 \in D$ and $\delta(1)=-1$.

Lemma 4. We have $-1 \in D$ if and only if $\varphi=\rho e^{\eta}$, where $\rho$ and $\eta$ are real-valued and $\eta(0)=0$ and $\eta$ is constant on each open interval on which $\rho$ (or $\varphi)$ does not vanish. Moreover, $\delta(-1)=1$.

Proof. Let $(x, y)$ be a maximal interval on which $\varphi$ does not vanish, so $\varphi$ does not vanish either on $(-y,-x)$ (we may have $(x, y)=\mathbb{R}$, of course). We can write $\varphi=e^{\psi}$ with $\psi$ of class $C^{1}$ on $(x, y)$ and $(-y,-x)$, and since $\psi(-v)=\overline{\psi(v)}$ the property $-1 \in D$ and (12) yield

$$
\psi^{\prime}(v)=\gamma^{\prime} \overline{\psi^{\prime}(v)} \quad \forall v \in(x, y) .
$$

Since $\gamma^{\prime} \in \mathbb{R}^{*}$, we deduce that $\psi^{\prime}(v) \in \mathbb{R}$ and thus $\gamma^{\prime}=1$ (because $\psi^{\prime}$ cannot be identically 0 ). Therefore, if $v_{0} \in(x, y)$, we have $\psi(v)-\psi\left(v_{0}\right) \in \mathbb{R}$ for all $v \in(x, y)$ and hence $\varphi=\rho e^{\eta}$ with $\eta(v)=\eta\left(v_{0}\right) \in \mathbb{R}$ for all $v \in(x, y)$. The converse is obvious.

Lemma 5. Let $\gamma \in \mathbb{R}^{*}$ with $|\gamma| \neq 1$. Then $\gamma \in D$ if and only if $\varphi$ does not vanish, and satisfies for some $C(\gamma)>0$

$$
\begin{equation*}
\varphi(v)=\varphi(\gamma v)^{C(\gamma)} \quad \forall v \in \mathbb{R} . \tag{14}
\end{equation*}
$$

Moreover,
a) $\operatorname{Re} \psi(v)<0$ for all $v \in \mathbb{R}^{*}$.
b) $\delta(\gamma)=-\gamma C(\gamma)$.
c) For all $n \in \mathbb{Z}$ we have $\gamma^{n} \in D$ and $C\left(\gamma^{n}\right)=C(\gamma)^{n}$.
d) $-\gamma \in D$ if and only if $\varphi$ is real-valued, and then $C(-\gamma)=C(\gamma)$.

Proof. The sufficient condition is obvious, as well as (b).
Conversely, assume that $\gamma \in D$. Let $(-x, x)$ be the maximal interval on which $\varphi$ does not vanish. We have $\varphi=e^{\psi}$ with $\psi$ of class $C^{1}$ on $(-x, x)$. For simplicity we set $\psi_{r}=\operatorname{Re} \psi$, and we have $\psi_{r}(u) \rightarrow-\infty$ as $|u| \uparrow x$ if $x<\infty$. On $(-x, x)$, (12) yields $\psi^{\prime}(v)+\gamma^{\prime} \psi^{\prime}(\gamma v)=0$, so $\psi(v)+\frac{\gamma^{\prime}}{\gamma} \psi(\gamma v)=0$, since $\psi(0)=0$.

If $|\gamma|>1$ and $x<\infty$, then $\left|\psi_{r}(v)\right|=\left|\frac{\gamma^{\prime}}{\gamma}\right|\left|\psi_{r}(\gamma v)\right| \rightarrow \infty$ as $|v| \uparrow x /|\gamma|$, contradicting the fact that $\psi$ is continuous on $(-x, x)$. Similarly, if $|\gamma|>1$ and $x<\infty,\left|\psi_{r}(\gamma v)\right|=\left|\frac{\gamma}{\gamma^{\prime}} \| \psi_{r}(v)\right| \rightarrow \infty$ as $|v| \uparrow x$, bringing up the same contradiction; therefore $x=\infty$, and $\varphi$ does not vanish. It follows that $\varphi=e^{\psi}$ everywhere and, with $C^{\prime}(\gamma)=-\gamma^{\prime} / \gamma$,

$$
\begin{equation*}
\psi(v)=C(\gamma) \psi(\gamma v) \quad \forall v \in \mathbb{R}, \tag{15}
\end{equation*}
$$

that is, we have (14). Since $U_{1}$ is non-degenerate, $\psi$ is not identically 0 and thus $C(\gamma) \neq 0$. Note also that (c) is obvious from (14).

We always have that $\psi_{r} \leq 0$ and that $\psi_{r}$ is even. Assume that $\psi_{r}(v)=0$ for some $v>0$. Then (15) and (c) imply $\psi_{r}\left(v|\gamma|^{n}\right)=0$ for all $n \in \mathbb{Z}$. It follows that the characteristic fonction of the symmetrized random variable $U=U_{1}-U_{2}$ equals 1 for all $v|\gamma|^{n}, n \in \mathbb{Z}$, so $U$ is supported by $\left\{2 k \pi / v|\gamma|^{n}: k \in \mathbb{Z}\right\}$, for all $n \in \mathbb{Z}$, which implies that $U=0$ a.s., contradicting again the non-degeneracy assumption. Thus (a) holds and (15) yields $C(\gamma)>0$.

Finally, it only remains to prove (d). If $\varphi$ is real-valued, it is even and (14) is satisfied with $-\gamma$ and $C(-\gamma)=C(\gamma)$. Suppose conversely that $-\gamma \in D$, then (15) gives $\overline{\psi(v)}=C(\gamma) \psi(-\gamma v)$, while $-\gamma \in D$ yields $\psi(v)=C(-\gamma) \psi(-\gamma v)$. Comparing the real parts of these two equalities and using (a) we obtain $C(-\gamma)=C(\gamma)$. Then $\bar{\psi}=\psi$ and $\varphi$ is real-valued.

Lemma 6. With $D_{+}=D \cap \mathbb{R}_{+}$, one of the following three cases is satisfied:
$\left(\mathbf{C}_{+} \mathbf{1}\right) \quad D_{+}=\{1\}$.
(C+2) $D_{+}=\left\{r^{n}: n \in \mathbb{Z}\right\}$ for some $r>1$.
(C+3) $\quad D_{+}=\mathbb{R}_{+}^{*}$.
Moreover, we are in case $\left(C_{+}\right.$3) if and only if either $\varphi(u)=\epsilon^{-a|u|^{2}}$ for some $a>0$ or $\varphi(u)=e^{-a|u|^{\rho}(1+i b s i g n(u))}$ for some $a>0, \rho \in(1,2),|b| \leq \tan \left(\frac{\pi}{2(2-\rho)}\right)$.

Proof. Due to the fact that $1 \in D$ and to Lemma 5, if we are not in case $\left(\mathrm{C}_{+} 1\right)$, $D_{+}$contains at least a $\gamma>0, \gamma \neq 1$, and then $\varphi=\epsilon^{\psi}$ satisfies (14). Indeed, $D_{+}$is the set of all $\gamma>0$ such that (15) holds for some $C(\gamma)>0$. Then $D_{+}$is clearly a multiplicative group, therefore it is closed since $\psi$ is continuous and thus it is of the form $\left(\mathrm{C}_{+} 2\right)$ or $\left(\mathrm{C}_{+} 3\right)$.

Assuming ( $\mathrm{C}_{+} 3$ ), for each $\gamma>0$ there is $C(\gamma)>0$ such that, if $f$ denotes either the real or the imaginary part of $\psi$, we have $f(0)=0$ and

$$
f(v)=C(\gamma) f(\gamma v) \quad \forall v \geq 0 .
$$

Then $f$ is either identically 0 , or everywhere positive, or everywhere negative, on $(0, \infty)$. In the last two cases, $g(u)=\log \left|f\left(e^{u}\right) / f(1)\right|$ satisfies $g(u+\log \gamma)=$ $g(u)+g(\log \gamma)$ for all $u \in \mathbb{R}, \gamma>0$, i.e., $g\left(u+u^{\prime}\right)=g(u)+g\left(u^{\prime}\right)$ for all $u, u^{\prime} \in \mathbb{R}$. Since $g$ is continuous, we obtain $g(u)=K u$. Thus, in all cases we have $f(v)=\eta v^{\rho}$ for some $\eta, \rho \in \mathbb{R}$, and furthermore $\gamma^{\rho} C(\gamma)=1$ for all $\gamma>0$ (hence $\rho$ is the same for both the real and imaginary parts of $\psi)$. We then deduce that $\psi(v)=(\alpha+i \beta) v^{\rho}$ for some $\alpha, \beta, \rho \in \mathbb{R}$, if $v>0$. By (a) of Lemma 5 we have $\alpha<0$ and since $\psi(-v)=\overline{\psi(v)}$, we also have $\psi(v)=(\alpha-i \beta)|v|^{\rho}$ for $v<0$. Then $\psi(v)=-a|v|^{\rho}(1+i b s i g n(v))$ for $a>0, b \in \mathbb{R}, \rho \in \mathbb{R}$. Conversely, each such $\psi$ satisfies (15) for all $\gamma>0$, with $C(\gamma 1)=\gamma^{-\rho}$, implying $D_{+}=\mathbb{R}_{+}^{*}$.

It remains to examine under which conditions on $(a, b, \rho)$ the function $\varphi=\epsilon^{\psi}$ with $\psi$ as above is a characteristic function. Observe that for all $\alpha, \alpha^{\prime}>0$ we have $\psi(\alpha v)+\psi\left(\alpha^{\prime} v\right)=\psi\left(\alpha^{\prime \prime} v\right)$ with $\alpha^{\prime \prime \rho}=\alpha^{\rho}+\alpha^{\prime \rho}$. Then, if it is the case, the corresponding distribution will be strictly stable, with a first moment equal to 0 . As is well known,
this will be the case if and only if either $\rho=2$ and $b=0$ (normal case), or $\rho \in(1,2)$ and $|b| \leq \tan \left(\frac{\pi}{2(2-\rho)}\right)$.

Proof of Theorem 3. a) When $U_{1}$ is symmetric, so is $D$, and $(\mathrm{Cs}-\mathrm{i})=\left(\mathrm{C}_{+} \mathrm{i}\right)$. Therefore Lemma 5 yields that one of (Cs-1), (Cs-2) or (Cs-3) is satisfied. Moreover, (Cs-2) implies that $\varphi$ never vanishes (by Lemma 5), and (Cs-3) holds if and only if $\varphi(v)=e^{-a|v|^{\rho}}$ (because here $\varphi$ is real-valued).
b) Now we suppose that $U_{1}$ is not symmetric. It suffices to prove that if $D \neq\{1\}$, then we are in one of the cases (Ca-i) for $\mathrm{i}=2,3,4,5$.

First, by Lemma $4,-1 \in D$ if and only if the necessary and sufficient condition in (Ca-2) is satisfied. Then $\varphi$ vanishes somewhere, and $D$ contains no $\gamma$ with $|\gamma| \neq 1$ by Lemma 5 . Thus $-1 \in D$ if and only if (Ca-2) holds.

Next, suppose that we are not in any of the cases (Ca-1) and (Ca-2). If $D=D_{+}$, we are then in cases (Ca-3) or (Ca-5) by Lemma 5. Otherwise there exists $\gamma>0$ with $\gamma \neq 1$ and $-\gamma \in D$. Then $\gamma^{2} \in D$ and $\gamma^{2} \neq 1$ and by Lemma 5 either ( $\mathrm{C}_{+} 2$ ) or ( $\mathrm{C}_{+} 3$ ) holds. However, under $\left(\mathrm{C}_{+} 3\right.$ ) we also have $\gamma \in D$, hence Lemma $5(\mathrm{~d})$ contradicts the assumption that $U_{1}$ is non-symmetric and indeed we have $\left(\mathrm{C}_{+} 2\right)$ with some $r>1$. It then follows that $\gamma^{2}=r^{k}$ for some $k \in N^{*}$, while Lemma 5(c) gives $C\left(r^{n}\right)=C(r)^{n}$ and $C(\gamma)=C(r)^{k / 2}$. Furthermore if $k$ were even we would have $r^{k / 2} \in D$ and $-r^{k / 2}=-\gamma \in D$, again a contradiction by Lemma $5(\mathrm{~d})$, so $k=2 p+1$ with $p \in \mathbb{Z}$ and $\gamma=r^{p+1 / 2}$. In order to obtain (Ca-4), it thus remains to prove that $-r^{n+1 / 2} \in D$ for all $n \in \mathbb{Z}$. For this, a repeated use of (15) yields

$$
\psi(v)=C\left(r^{n-p}\right) \psi\left(r^{n-p} v\right)=C\left(r^{n-p}\right) C(\gamma) \psi\left(-\gamma r^{n-p} v\right)=C(r)^{n+1 / 2} \psi\left(-r^{n+1 / 2}\right)
$$

and the result follows.
c) Since $\delta(1)=-1$ and $\delta(-1)=1$ (Lemma 4), (c) is obvious in cases (Cs-1), (Ca1) and (Ca-2). Also, (c) with $\theta=\rho$ follows from Lemma $5(\mathrm{~b})$ and from a comparison between (14) and the explicit form of $\varphi$ in cases (Cs-3) and (Ca-5).

Under (Ca-4) we have seen that $C^{C}\left(r^{n}\right)=C^{n}$ and $C\left(-r^{n+1 / 2}\right)=C^{n+1 / 2}$, where $C=$ $C(r)$. Thus, for all $\gamma \in D, \quad C(\gamma)=C^{\log (|\gamma|) / \log (r)}=|\gamma|^{-\theta} \quad$ with $\theta=-\log (C) / \log (r)$ (so $\left.\delta(\gamma)=-\gamma /|\gamma|^{\theta}\right)$. The same holds for (Cs-2) and (Ca-3). Now (15) yields $\psi^{\prime}(v)=$ $C r \psi^{\prime}(r v)$, hence $\psi^{\prime}\left(r^{-n} v\right)=(C r)^{n} \psi^{\prime}(v)$ and since $\psi^{\prime}(0)=0$ and $\psi^{\prime}(v) \neq 0$ for some $v \neq 0$ we must have $C r<1$ and thus $\theta>1$.

Suppose that $\theta>2$, i.e. $A:=C r^{2}<1$. Then $\psi^{\prime}\left(r^{-n} v\right) /\left(r^{-n} v\right)=A^{n} \psi^{\prime}(v) / v$ and if $|w| \leq r^{-n}$ there is $m \geq n$ and $v \in(1 / r, 1]$ with $w=v r^{-m}$ or $w=-v r^{-m}$. Therefore $\sup _{|w| \leq r^{-n}}\left|\psi^{\prime}(w) / w\right| \leq A^{n} \sup _{1 / r<|v| \leq 1}\left|\psi^{\prime}(v)\right|$. It follows that $\psi^{\prime}$ is differentiable at 0 , with $\bar{\psi}^{\prime \prime}(0)=0$. Hence $U_{1}$ is square-integrable, with variance 0 , which contradicts once more the non-degeneracy assumption and therefore $\theta \leq 2$.
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