# CHARACTERISTIC FUNCTIONS OF RANDOM VARIABLES ATTRACTED TO 1-STABLE LAWS ${ }^{1}$ 

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#### Abstract

The domain of attraction of a 1-stable law on $\mathbb{R}^{d}$ is characterized by the expansions of the characteristic functions of its elements.


0 . Introduction. Let $X_{1}, X_{2}, \ldots$ be $\mathbb{R}^{d}$-valued, independent, identically distributed random variables. The distributional limits of $\left(S_{n}-A_{n}\right) / B_{n}$, where $A_{n} \in \mathbb{R}^{d}, B_{n}>0$ are constants and $S_{n}=\sum_{k=1}^{n} X_{k}$, are given by the well-known stable laws. [Lévy (1954), Gnedenko and Kolmogorov (1954) and I bragimov and Linnik (1971)].

A probability distribution function $F$ on $\mathbb{R}^{d}$ is called stable if for all $a, b>0$ there are $c>0$ and $v \in \mathbb{R}^{d}$ such that

$$
F_{a} * F_{b}(x)=F_{c}(x-v), \quad x \in \mathbb{R}^{d},
$$

where $F_{s}(x)=F(x / s), x \in \mathbb{R}^{d}, s>0$, and strictly stableif this is true with $v=0$.
In this case [Lévy (1954)] necessarily $a^{p}+b^{p}=c^{p}$ for some $0<p \leq 2$, and $p$ is called the order of the stable law $F$.

A distribution $G$ on $\mathbb{R}^{d}$ belongs to the domain of attraction of the stable law $F$ if there are constants $A_{n} \in \mathbb{R}^{d}$ and $B_{n}>0$ such that the distributions $\left(S_{n}-A_{n}\right) / B_{n}$ converge weakly to $F$ where $S_{n}=X_{1}+\cdots+X_{n}$ and $X_{1}, X_{2}, \ldots$ are i.i.d. with distribution $G$.

For $p \in(0,2]$ and $d \in \mathbb{N}$, we let $\operatorname{DA}(p, d)$ be the collection of distribution functions in the domain of attraction of some stable law on $\mathbb{R}^{d}$ of order $p$.

In this paper, we obtain expansions of the characteristic functions of distributions on $\mathbb{R}^{d}$ which are in the domain of attraction of a stable law.

In Section 1 we deal with the case $d=1$. The first partial results are in Gnedenko and Koroluk (1950). The expansions are given fully in Ibragimov and Linnik (1971) in case $p \neq 1$ (see Theorem 1).

Our main result is Theorem 2 giving the expansions in case $p=1$.
In Section 2 we obtain as corollaries expansions in case $d \geq 2$. Other results in this case are to be found in Rvac̃eva (1962), Meerschaert (1986), Kuelbs and Mandrekar (1974) and Araujo and Giné (1979, 1980).

A stable law of order $p$ on $\mathbb{R}$ has a characteristic function $\psi$ of the form

$$
\log \psi(t)=i t \gamma-c|t|^{p}\left[1-i \beta \operatorname{sgn}(t) \tan \left(\frac{p \pi}{2}\right)\right], \quad p \neq 1,
$$

[^0]and
$$
\operatorname{Relog} \psi(t)=-c|t|, \quad \operatorname{Im} \log \psi(t)=t\left(\gamma+\frac{2 \beta c}{\pi} \log \left(\frac{1}{|t|}\right)\right), \quad p=1
$$
where $c>0, \beta, \gamma \in \mathbb{R}$ are constants [Lévy (1954)].
The form of the characteristic functions of stable laws on $\mathbb{R}^{d}$ was obtained by Feldheim [see Feldheim (1937), Lévy (1954) and Samorodnitsky and Taqqu (1994), Theorem 2.3.1]:

To each stable law of order $p$ on $\mathbb{R}^{d}$ there corresponds a finite measure $\nu$ on $S^{d-1}$ (called the spectral measure) and $\mu \in \mathbb{R}^{d}$ (called the translate) so that the characteristic function $\psi$ has the form

$$
\begin{equation*}
\log \psi(u)=i\langle u, \mu\rangle-\int_{S^{d-1}}|\langle u, s\rangle|^{p}\left(1-i \operatorname{sgn}(\langle s, u\rangle) \tan \left(\frac{p \pi}{2}\right)\right) \nu(d s) \tag{1a}
\end{equation*}
$$

for $p \neq 1$ and
(1b) $\quad \log \psi(u)=i\langle u, \mu\rangle-\int_{S^{d-1}}|\langle u, s\rangle|\left(1+i \frac{2}{\pi} \operatorname{sgn}(\langle u, s\rangle) \log (|\langle u, s\rangle|)\right) \nu(d s)$
for $p=1$. Evidently a stable law on $\mathbb{R}^{d}$ has a density if and only if the support of its spectral measure is not contained in a proper subspace of $\mathbb{R}^{d}$, and in this case we say that both the stable law and the spectral measure are nondegenerate.

Clearly, the stability of an $\mathbb{R}^{d}$-valued random variable $Z$ implies that of its inner products $\langle Z, u\rangle, u \in \mathbb{R}^{d}$.

An example of Marcus (1983) shows that the converse of this is false without additional assumptions.

According to Theorems 2.1.2 and 2.1.5 in Samorodnitsky and Taqqu (1994), the $\mathbb{R}^{d}$-valued random variable $Z$ is strictly stable (stable with index $\geq 1$ ) if its inner products $\langle Z, u\rangle, u \in \mathbb{R}^{d}$, are strictly stable on $\mathbb{R}$ (stable on $\mathbb{R}$ with index $\geq 1$ ).

The first characterizations of domains of attraction were in terms of the tails of the distributions concerned.

In the unidimensional case [Gnedenko and Kolmogorov (1954)], for $p<2$, the (right continuous) distribution function $G \in \mathrm{DA}(p, 1)$ iff there is a function $L: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$, slowly varying at $\infty$ [see Feller (1971)], and constants $c_{1}, c_{2} \geq 0$, $c_{1}+c_{2}>0$ such that

$$
\begin{align*}
& L_{1}(x):=x^{p}(1-G(x))=\left(c_{1}+o(1)\right) L(x),  \tag{2}\\
& L_{2}(x):=x^{p} G(-x)=\left(c_{2}+o(1)\right) L(x) \quad \text { as } x \rightarrow+\infty .
\end{align*}
$$

The results of Gnedenko and Kolmogorov (1954) were generalized to $\mathbb{R}^{d}$ in Rvac̃eva (1962) [see also M eerschaert (1986)], to Hilbert space in Kuelbs and Mandrekar (1974), and to Banach space in Araujo and Giné (1979).

1. Unidimensional characterization. The characteristic function $\psi$ of $G \in \mathrm{DA}(p, 1)$ is considered in Gnedenko and Koroluk (1950) and Ibragimov and Linnik (1971).

In Gnedenko and K oroluk (1950), $\mathrm{DA}(p, 1)$ is characterized in terms of $\psi(t)$.
In Ibragimov and Linnik (1971), the asymptotic expansion of $\log \psi(t)$ around 0 is established with error small when compared to

$$
\operatorname{Prob} .\left(|Z|>\frac{1}{|t|}\right)=|t|^{p}\left(L_{1}\left(\frac{1}{|t|}\right)+L_{2}\left(\frac{1}{|t|}\right)+\right)=|t|^{p}\left(c_{1}+c_{2}+o(1)\right) L\left(\frac{1}{|t|}\right)
$$

as $t \rightarrow 0$. Here, $Z$ is a $G$-distributed random variable, and $G \in \operatorname{DA}(p, 1)$, $p \neq 1$, satisfies (2) with the slowly varying functions $L, L_{1}, L_{2}$ and constants $c_{1}, c_{2} \geq 0, c_{1}+c_{2}>0$. Specifically:

Theorem 1 [Ibragimov and Linnik (1971), Theorem 2.6.5]. Suppose that $G$ satisfies (2) with $p \neq 1$. Then

$$
\log \psi(t)=i t \gamma-c|t|^{p} L\left(|t|^{-1}\right)\left[1-i \beta \operatorname{sgn}(t) \tan \left(\frac{p \pi}{2}\right)\right]+o\left(|t|^{p} L\left(|t|^{-1}\right)\right)
$$

where

$$
\begin{gathered}
\beta=\frac{c_{1}-c_{2}}{c_{1}+c_{2}}, \quad c=\Gamma(1-p)\left(c_{1}+c_{2}\right) \cos \left(\frac{p \pi}{2}\right), \\
\gamma= \begin{cases}0, & 0<p<1 \\
\int x G(d x), & 1<p \leq 2\end{cases}
\end{gathered}
$$

The expansion of the characteristic function when $p=1$ is also treated in I bragimov and Linnik (1971) for a limited class of slowly varying functions $L$, namely those where

$$
\int_{0}^{\lambda} \frac{x L(x) d x}{1+x^{2}}=L(\lambda)(\log \lambda+o(1))
$$

as $\lambda \rightarrow \infty$ [cf. Theorem 2 here, Theorem 2.6 .5 there and formula (2.6.34) there]. As can be easily checked, the functions $L(x) \sim(\log x)^{a}, a \in \mathbb{R}$, and $L(x) \sim \exp \left[(\log x)^{a}\right], 0<a<1$, are slowly varying functions not in this class.

Theorem 2. Suppose that $G$ satisfies (2) with $p=1$. Then

$$
\begin{aligned}
& \operatorname{Relog} \psi(t)=-c|t| L\left(|t|^{-1}\right)+o\left(|t| L\left(|t|^{-1}\right)\right) \\
& \operatorname{Im} \log \psi(t)=t \gamma+\frac{2 \beta c}{\pi} C t L\left(\frac{1}{|t|}\right)+t\left(H_{1}\left(\frac{1}{|t|}\right)-H_{2}\left(\frac{1}{|t|}\right)\right)+o\left(|t| L\left(|t|^{-1}\right)\right)
\end{aligned}
$$

as $t \rightarrow 0$, where

$$
\begin{aligned}
H_{j}(\lambda) & =\int_{0}^{\lambda} \frac{x L_{j}(x) d x}{1+x^{2}}, \quad j=1,2, \\
C & =\int_{0}^{\infty}\left(\cos y-\frac{1}{1+y^{2}}\right) \frac{d y}{y},
\end{aligned}
$$

and the constants $c>0, \beta, \gamma \in \mathbb{R}$ are defined by

$$
\begin{aligned}
\beta & =\frac{c_{1}-c_{2}}{c_{1}+c_{2}}, \quad c=\frac{\left(c_{1}+c_{2}\right) \pi}{2} \\
\gamma & =\int_{-\infty}^{\infty}\left(\frac{x}{1+x^{2}}+\operatorname{sgn}(x) \int_{0}^{|x|} \frac{2 u^{2}}{\left(1+u^{2}\right)^{2}} d u\right) G(d x)
\end{aligned}
$$

Remark 1. Note that

$$
H_{1}(\lambda)=\int_{0}^{\lambda} \frac{x^{2} P(Z>x) d x}{1+x^{2}}
$$

whence

$$
\begin{aligned}
H_{1}(\lambda)-H_{2}(\lambda) & =E\left(\left[|Z| \wedge \lambda-\tan ^{-1}(|Z| \wedge \lambda)\right] \operatorname{sgn}(Z)\right) \\
& =E((|Z| \wedge \lambda) \operatorname{sgn}(Z))+O(1)
\end{aligned}
$$

as $\lambda \rightarrow \infty$, where $Z$ is $G$-distributed and $H_{1}, H_{2}$ are as in Theorem 2.
REMARK 2. From this representation of the characteristic function of distributions in $\operatorname{DA}(p, 1)$, one deduces the existence of a $p$-stable random variable $Y$ and constants $A_{n}, B_{n} \in \mathbb{R}, B_{n}>0$ so that $\left(S_{n}-A_{n}\right) / B_{n} \rightarrow Y$ in distribution. These constants [unique up to $o\left(B_{n}\right)$ as $n \rightarrow \infty$ ] are given by

$$
n L\left(B_{n}\right)=B_{n}^{p}, \quad A_{n}= \begin{cases}0, & 0<p<1 \\ \gamma n, & 1<p \leq 2 \\ \gamma n+n\left(H_{1}\left(B_{n}\right)-H_{2}\left(B_{n}\right)\right), & p=1\end{cases}
$$

To see this in case $p=1$, write

$$
\log E\left(\exp \left[i t\left(\frac{S_{n}-A_{n}}{B_{n}}\right)\right]\right)=-\frac{i t A_{n}}{B_{n}}+n \log \psi\left(\frac{t}{B_{n}}\right):=\alpha_{n}(t)+i \beta_{n}(t)
$$

Then

$$
\alpha_{n}(t)=-c \frac{n|t|}{B_{n}} L\left(\frac{B_{n}}{|t|}\right)+o\left(\frac{n|t| L\left(B_{n} /|t|\right)}{B_{n}}\right) \rightarrow-c|t| \quad \text { as } n \rightarrow \infty
$$

and

$$
\begin{aligned}
\beta_{n}(t)= & \frac{t\left(H_{1}\left(B_{n} /|t|\right)-H_{1}\left(B_{n}\right)\right)}{L\left(B_{n}\right)}-\frac{t\left(H_{2}\left(B_{n} /|t|\right)-H_{2}\left(B_{n}\right)\right)}{L\left(B_{n}\right)} \\
& +\frac{2 \beta c t C L\left(B_{n} /|t|\right)}{\pi L\left(B_{n}\right)}+o\left(\frac{n|t| L\left(B_{n} /|t|\right)}{B_{n}}\right)
\end{aligned}
$$

Now, for $j=1,2$ and $k>1$ [see (5) in Lemma 3 below],

$$
H_{j}(k \lambda)-H_{j}(\lambda)=c_{j} L(\lambda) \log k+o(L(\lambda)) \quad \text { as } \lambda \rightarrow \infty
$$

Thus, with $k=1 /|t|$,

$$
\beta_{n}(t) \rightarrow t\left(c_{1}-c_{2}\right) \log \frac{1}{|t|}+\frac{2 \beta c C t}{\pi}=\frac{2 \beta c t}{\pi}\left(\log \frac{1}{|t|}+C\right) \quad \text { as } n \rightarrow \infty .
$$

Thus, the above representation is a characterization of $\operatorname{DA}(p, 1)$.
Remark 3. We notethat the expansion of $\psi(t)$ around 0 up to $o\left(|t|^{p} L(1 /|t|)\right)$ is determined entirely by the asymptotic equivalence class of the slowly varying function $L$ and the constants $c_{1}, c_{2} \geq 0$ for $G$ satisfying (2) with $p \neq 1$.

This is not the case when $p=1$ as shown by the following examples.
There is a distribution $G$ so that

$$
\begin{aligned}
& L_{1}(x):=x(1-G(x))=(\log x)^{2}+(\log x)^{3 / 2}+O(1), \\
& L_{2}(x):=x G(-x)=(\log x)^{2}+O(1) \quad \text { as } x \rightarrow+\infty .
\end{aligned}
$$

Here, $L(\lambda)=(\log \lambda)^{2}, p=c_{1}=c_{2}=1$, and one calculates from Theorem 2 that

$$
\operatorname{Im} \log \psi(t)=\frac{4 t}{5 \pi} L\left(\frac{1}{|t|}\right)^{5 / 4}+o\left(|t| L\left(\frac{1}{|t|}\right)\right) \quad \text { as } t \rightarrow 0
$$

On the other hand, there is a symmetric distribution satisfying

$$
L_{1}(x)=L_{2}(x)=(\log x)^{2}+O(1) \quad \text { as } x \rightarrow+\infty
$$

for which also $L(\lambda)=(\log \lambda)^{2}$, and $p=c_{1}=c_{2}=1$; but here (owing to symmetry)

$$
\operatorname{Im} \log \psi(t) \equiv 0
$$

Proof of Theorem 2. Assume that $G$ is represented in the form (2).
For $x>0$ define distribution functions $G_{j}, j=1,2$, on $\mathbb{R}_{+}$by

$$
G_{1}(x)=G(x)-G(0) \quad \text { and } \quad G_{2}(x)=G(0)-G(-x) .
$$

We have that

$$
G_{j}(\infty)-G_{j}(x)=\frac{L_{j}(x)}{x}=\frac{\left(c_{j}+o(1)\right) L(x)}{x} .
$$

Write

$$
\begin{aligned}
\int\left(1-\exp (i t x)+\frac{i t x}{1+x^{2}}\right) G(d x)= & \int_{0}^{\infty}\left(1-\exp (i t x)+\frac{i t x}{1+x^{2}}\right) G_{1}(d x) \\
& +\int_{0}^{\infty}\left(1-\frac{i t x}{1+x^{2}}-\exp (-i t x)\right) G_{2}(d x)
\end{aligned}
$$

and let

$$
\gamma_{j}=\int_{0}^{\infty} \frac{2 x^{2}}{\left(1+x^{2}\right)^{2}}\left(G_{j}(\infty)-G_{j}(x)\right) d x=\int_{0}^{\infty} \frac{2 x L_{j}(x) d x}{\left(1+x^{2}\right)^{2}} .
$$

Integration by parts gives

$$
\begin{aligned}
\int_{0}^{\infty} & \left(1-\exp \left[-(-1)^{j} i t x\right]-(-1)^{j} \frac{i t x}{1+x^{2}}\right) G_{j}(d x) \\
& =(-1)^{j} i t \int_{0}^{\infty}\left(\exp \left[-(-1)^{j} i t x\right]-\frac{1-x^{2}}{\left(1+x^{2}\right)^{2}}\right) \frac{L_{j}(x) d x}{x} \\
& =|t| \int_{0}^{\infty} \sin (|t| x) \frac{L_{j}(x) d x}{x}+(-1)^{j} i t \int_{0}^{\infty}\left(\cos (t x)-\frac{1-x^{2}}{\left(1+x^{2}\right)^{2}}\right) \frac{L_{j}(x) d x}{x}
\end{aligned}
$$

Changing variables, we obtain that

$$
\begin{gathered}
\int_{0}^{\infty} \sin (|t| x) \frac{L_{j}(x) d x}{x}=\int_{0}^{\infty} \sin (x) \frac{L_{j}(x /|t|) d x}{x} \\
\int_{0}^{\infty}\left(\cos (t x)-\frac{1}{1+(t x)^{2}}\right) \frac{L_{j}(x) d x}{x}=\int_{0}^{\infty}\left(\cos (x)-\frac{1}{1+x^{2}}\right) \frac{L_{j}(x /|t|) d x}{x} .
\end{gathered}
$$

By Lemma 1, we see that

$$
\int_{0}^{\infty} \sin (|t| x) \frac{L_{j}(x) d x}{x}=(1+o(1)) L_{j}\left(\frac{1}{|t|}\right) \frac{\pi}{2}
$$

Now

$$
\begin{aligned}
\int_{0}^{\infty}( & \left.\cos (t x)-\frac{1-x^{2}}{\left(1+x^{2}\right)^{2}}\right) \frac{L_{j}(x) d x}{x} \\
= & \int_{0}^{\infty}\left(\cos (t x)-\frac{1}{1+(t x)^{2}}\right) \frac{L_{j}(x) d x}{x}+\int_{0}^{\infty} \frac{x\left(1-t^{2}\right) L_{j}(x) d x}{\left(1+x^{2}\right)\left(1+(t x)^{2}\right)} \\
& +\int_{0}^{\infty} \frac{2 x L_{j}(x) d x}{\left(1+x^{2}\right)^{2}} \\
= & \int_{0}^{\infty}\left(\cos (t x)-\frac{1}{1+(t x)^{2}}\right) \frac{L_{j}(x) d x}{x}+\int_{0}^{\infty} \frac{x\left(1-t^{2}\right) L_{j}(x) d x}{\left(1+x^{2}\right)\left(1+(t x)^{2}\right)}+\gamma_{j}
\end{aligned}
$$

By Lemma 2,

$$
\int_{0}^{\infty}\left(\cos (t x)-\frac{1}{1+(t x)^{2}}\right) \frac{L_{j}(x) d x}{x}=C L_{j}\left(\frac{1}{|t|}\right)+o\left(L\left(\frac{1}{|t|}\right)\right)
$$

Set

$$
\tilde{H}_{j}(\lambda):=\int_{0}^{\infty} \frac{x L_{j}(x) d x}{\left(1+x^{2}\right)\left(1+x^{2} / \lambda^{2}\right)}
$$

By Lemma 3, $\tilde{H}_{j}(\lambda)=H_{j}(\lambda)+o(L(\lambda))$ as $\lambda \rightarrow \infty$.

Putting everything together, we obtain

$$
\begin{aligned}
\int_{0}^{\infty}(1+ & \left.\frac{i t x}{1+x^{2}}-\exp (i t x)\right) G_{1}(d x)+\int_{0}^{\infty}\left(1-\frac{i t x}{1+x^{2}}-\exp (-i t x)\right) G_{2}(d x) \\
= & L\left(\frac{1}{|t|}\right)|t|\left(c_{1}+c_{2}\right) \frac{\pi}{2}-i t L\left(\frac{1}{|t|}\right)\left(c_{1}-c_{2}\right) C \\
& -i t\left(\tilde{H}_{1}\left(\frac{1}{|t|}\right)-\tilde{H}_{2}\left(\frac{1}{|t|}\right)\right)-i t\left(\gamma_{1}-\gamma_{2}\right)+o\left(|t| L\left(\frac{1}{|t|}\right)\right) \\
= & L\left(\frac{1}{|t|}\right)|t|\left(c_{1}+c_{2}\right) \frac{\pi}{2}-i t L\left(\frac{1}{|t|}\right)\left(c_{1}-c_{2}\right) C \\
& -i t\left(H_{1}\left(\frac{1}{|t|}\right)-H_{2}\left(\frac{1}{|t|}\right)\right)-i t\left(\gamma_{1}-\gamma_{2}\right)+o\left(|t| L\left(\frac{1}{|t|}\right)\right)
\end{aligned}
$$

and hence Theorem 2.
We conclude this section by collecting the lemmas on slowly varying functions needed for Theorem 2.

Assume that $h: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$is locally integrable, slowly varying at infinity and such that $u \mapsto h(u) / u$ is a nonincreasing function. Recall that $h$ has a representation

$$
h(x)=\eta(x) \exp \left[\int_{1}^{x} \frac{\varepsilon(s)}{s} d s\right]
$$

for some functions $\eta(s) \rightarrow K \in \mathbb{R}$ and $\varepsilon(s) \rightarrow 0$ as $s \rightarrow \infty$ [see Feller (1971)].
Lemma 1.

$$
\int_{0}^{\infty} \frac{\sin y}{y} h\left(\frac{y}{t}\right) d y=(1+o(1)) h\left(\frac{1}{t}\right) \frac{\pi}{2} .
$$

Proof. As the proof of Lemma 2.6.1 in Ibragimov and Linnik (1971).
Lemma 2.

$$
\int_{0}^{\infty}\left[\cos y-\frac{1}{1+y^{2}}\right] \frac{1}{y} h\left(\frac{y}{t}\right) d y=(1+o(1)) h\left(\frac{1}{t}\right) \int_{0}^{\infty}\left[\cos y-\frac{1}{1+y^{2}}\right] \frac{1}{y} d y .
$$

Proof. We first split the region of integration into four parts: $I_{1}=\left[\Delta_{1}, \infty\right)$, $I_{2}=\left[\delta, \Delta_{1}\right), I_{3}=\left[t \Delta_{2}, \delta\right)$ and $I_{4}=\left[0, t \Delta_{2}\right)$ where $\delta<1<\Delta_{1}=\left(N-\frac{1}{2}\right) \pi$, $N \in \mathbb{N}$.

Since $\left|\int_{\left[\Delta_{1}+n \pi, \Delta_{1}+(n+1) \pi\right]} \cos y[h(y / t) d y / y]\right|$ decreases in $n$,

$$
\left|\int_{I_{1}} \cos y \frac{h(y / t) d y}{y}\right| \leq \frac{\pi h\left(\Delta_{1} / t\right)}{\Delta_{1}} \sim \frac{\pi h(1 / t)}{\Delta_{1}} .
$$

Also,

$$
\int_{I_{1}} \frac{1}{1+y^{2}} \frac{h(y / t) d y}{y} \leq \frac{h\left(\Delta_{1} / t\right)}{\Delta_{1}} \pi \sim \frac{\pi h(1 / t)}{\Delta_{1}}
$$

Since, for $x \in\left[\Delta_{2} t, \delta\right)$,

$$
\frac{h(x / t)}{h(1 / t)}=(1+o(1)) \exp \left[\int_{x / t}^{1 / t} \frac{\varepsilon(s)}{s} d s\right]=\exp [o(-\log x)] \leq x^{-1 / 2}
$$

for $t$ small enough and $\Delta_{2}$ large enough,

$$
\begin{aligned}
\left|\int_{I_{3}}\left(\frac{1}{1+y^{2}}-\cos y\right) h\left(\frac{y}{t}\right) \frac{d y}{y}\right| & =O\left(h\left(\frac{1}{t}\right) \int_{0}^{\delta}\left|\frac{1}{1+y^{2}}-\cos y\right| y^{-3 / 2} d y\right) \\
& =O\left(h\left(\frac{1}{t}\right) \delta^{3 / 2}\right)
\end{aligned}
$$

Since the function $h$ is locally integrable, it follows that for $t$ small enough

$$
\begin{aligned}
\left|\int_{I_{4}}\left(\frac{1}{1+y^{2}}-\cos y\right) h\left(\frac{y}{t}\right) \frac{d y}{y}\right| & =\left|\int_{0}^{\Delta_{2}}\left(\frac{1}{1+t^{2} z^{2}}-\cos t z\right) h(z) \frac{d z}{z}\right| \\
& =O\left(t^{2} \Delta_{2} \int_{0}^{\Delta_{2}}|h(z)| d z\right) \\
& =O\left(t^{2}\right)=o\left(h\left(\frac{1}{t}\right)\right)
\end{aligned}
$$

For $\delta \leq x \leq \Delta_{1}$ we have (uniformly in $x$ ), by the slow variation property of $h$,

$$
\lim _{t \rightarrow 0} \frac{h(x / t)}{h(1 / t)}=1
$$

It follows that

$$
\begin{aligned}
& \left|\int_{I_{2}}\left(\frac{1}{1+y^{2}}-\cos y\right)\left[h\left(\frac{y}{t}\right)-h\left(\frac{1}{t}\right)\right] \frac{d y}{y}\right| \\
& \quad \leq 2 h\left(\frac{1}{t}\right)\left[\sup _{\delta \leq x \leq \Delta_{1}}\left|\frac{h(x / t)}{h(1 / t)}-1\right|\right] \int_{\delta}^{\Delta_{1}} \frac{d y}{y} \\
& \quad=o\left(h\left(\frac{1}{t}\right)\right)
\end{aligned}
$$

Applying the estimates for $I_{1}, I_{3}$ and $I_{4}$ with $h=1$, it follows that

$$
\int_{0}^{\infty}\left(\frac{1}{1+y^{2}}-\cos y\right) \frac{h(y / t)-h(1 / t)}{y} d y=o\left(h\left(\frac{1}{t}\right)\right)+O\left(h\left(\frac{1}{t}\right)\left(\delta^{3 / 2}+\Delta_{1}^{-1}\right)\right)
$$

Letting $\Delta_{1} \rightarrow \infty$ and $\delta \rightarrow 0$ as $t \rightarrow 0$, the lemma follows.

Lemma 3. Let

$$
H(\lambda):=\int_{0}^{\lambda} \frac{x h(x) d x}{1+x^{2}}
$$

then $H$ is slowly varying at $\infty$,

$$
\begin{equation*}
\frac{h(\lambda)}{H(\lambda)} \rightarrow 0 \quad \text { as } \lambda \rightarrow \infty \tag{3}
\end{equation*}
$$

$$
\begin{aligned}
\tilde{H}(\lambda) & :=\int_{0}^{\infty} \frac{x h(x) d x}{\left(1+x^{2}\right)\left(1+x^{2} / \lambda^{2}\right)} \\
& =H(\lambda)+o(h(\lambda)) \quad \text { as } \lambda \rightarrow \infty
\end{aligned}
$$

and

$$
\begin{equation*}
H(k \lambda)-H(\lambda) \sim h(\lambda) \cdot \log k \quad \text { as } \lambda \rightarrow \infty \tag{5}
\end{equation*}
$$

Remark. Slow variation of $H$, (3) and (5) are established in Lemma 1 of Parameswaran (1961).

Proof. We first show (5):

$$
\begin{aligned}
H(k \lambda)-H(\lambda) & =\int_{\lambda}^{k \lambda} \frac{x h(x) d x}{1+x^{2}} \sim \int_{\lambda}^{k \lambda} \frac{h(x) d x}{x} \\
& =\int_{1}^{k} \frac{h(\lambda x) d x}{x} \sim \log k h(\lambda) .
\end{aligned}
$$

Next, we see that (3) follows from (5) as $\forall M>1$,

$$
\begin{aligned}
\frac{H(\lambda)}{h(\lambda)} & =\frac{H\left(e^{M} e^{-M} \lambda\right)}{h(\lambda)} \\
& \geq \frac{H\left(e^{M} e^{-M} \lambda\right)-H\left(e^{-M} \lambda\right)}{h(\lambda)} \\
& \sim \frac{h\left(e^{-M} \lambda\right) M}{h(\lambda)} \rightarrow M \quad \text { as } \lambda \rightarrow \infty .
\end{aligned}
$$

It follows from (3) and (5) that $H$ is slowly varying at $\infty$.
To continue, we claim that
(6) $\quad \tilde{H}(\lambda)=\int_{0}^{\lambda} \frac{x h(x) d x}{\left(1+x^{2}\right)\left(1+x^{2} / \lambda^{2}\right)}+\frac{\log 2}{2} h(\lambda)+o(h(\lambda)) \quad$ as $\lambda \rightarrow \infty$.

To see this, note that

$$
\begin{aligned}
\int_{\lambda}^{\infty} \frac{x h(x) d x}{\left(1+x^{2}\right)\left(1+x^{2} / \lambda^{2}\right)}= & \int_{1}^{\infty} \frac{x h(\lambda x) d x}{\left(1 / \lambda^{2}+x^{2}\right)\left(1+x^{2}\right)} \\
= & h(\lambda) \int_{1}^{\infty} \frac{x d x}{\left(1 / \lambda^{2}+x^{2}\right)\left(1+x^{2}\right)} \\
& +h(\lambda) \int_{1}^{\infty}\left(\frac{h(\lambda x)}{h(\lambda)}-1\right) \frac{x d x}{\left(1 / \lambda^{2}+x^{2}\right)\left(1+x^{2}\right)} \\
= & \frac{\log 2}{2} h(\lambda)+o(h(\lambda))
\end{aligned}
$$

as $\lambda \rightarrow \infty$ by the dominated convergence theorem since $|h(\lambda x) / h(\lambda)-1| \rightarrow 0$ as $\lambda \rightarrow \infty \forall x>1$ and $|h(\lambda x) / h(\lambda)-1| \leq x \forall x$ large enough. This establishes (6).

To complete the proof of (4), we note that

$$
\frac{x h(x)}{\left(1+x^{2}\right)\left(1+x^{2} / \lambda^{2}\right)}=\frac{\lambda^{2}}{\lambda^{2}-1}\left(\frac{x h(x)}{x^{2}+1}-\frac{x h(x)}{x^{2}+\lambda^{2}}\right)
$$

whence, in view of (6),

$$
\tilde{H}(\lambda)=\frac{\lambda^{2}}{\lambda^{2}-1} \int_{0}^{\lambda} \frac{x h(x) d x}{x^{2}+1}-\frac{\lambda^{2}}{\lambda^{2}-1} \int_{0}^{\lambda} \frac{x h(x) d x}{x^{2}+\lambda^{2}}+\frac{\log 2}{2} h(\lambda)+o(h(\lambda))
$$

Now

$$
\begin{aligned}
\frac{\lambda^{2}}{\lambda^{2}-1} \int_{0}^{\lambda} \frac{x h(x) d x}{x^{2}+1} & =H(\lambda)+O\left(\frac{H(\lambda)}{\lambda^{2}}\right) \\
& =H(\lambda)+o(h(\lambda)) \quad \text { as } \lambda \rightarrow \infty
\end{aligned}
$$

because both $h$ and $H$ are slowly varying at $\infty$; and

$$
\begin{aligned}
\frac{\lambda^{2}}{\lambda^{2}-1} \int_{0}^{\lambda} \frac{x h(x) d x}{x^{2}+\lambda^{2}} & \sim \int_{0}^{\lambda} \frac{x h(x) d x}{x^{2}+\lambda^{2}} \\
& =\int_{0}^{1} \frac{x h(\lambda x) d x}{x^{2}+1} \\
& \sim \frac{\log 2}{2} h(\lambda) \text { as } \lambda \rightarrow \infty
\end{aligned}
$$

Thus,

$$
\tilde{H}(\lambda)=H(\lambda)+o(h(\lambda)) \quad \text { as } \lambda \rightarrow \infty
$$

which is (4).
2. Multidimensional characterization.

Corollary 1. Let $0<p<2, p \neq 1$ and $G$ be a distribution function on $\mathbb{R}^{d}$. The following are equi valent:
(A) $G$ belongs to the domain of attraction of the nondegenerate stable law of order $p$, spectral measure $\nu$ and translate $\mu$.
(B) The characteristic function $\psi$ of $G$ has the form

$$
\log \psi(t u)= \begin{cases}-t^{p} L\left(\frac{1}{t}\right) \Phi(u)+i t\langle u, \mu\rangle+o\left(t^{p} L\left(\frac{1}{t}\right)\right), & \text { if } p>1, \\ -t^{p} L\left(\frac{1}{t}\right) \Phi(u)+o\left(t^{p} L\left(\frac{1}{t}\right)\right), & \text { if } p<1\end{cases}
$$

as $t \rightarrow 0^{+}, \forall u \in S^{d-1}$, where $\mu \in \mathbb{R}^{d}, L$ is slowly varying at $\infty, \nu$ is a nondegenerate finite measure on $S^{d-1}$ and

$$
\Phi(u):=\int_{S^{d-1}}|\langle u, s\rangle|^{p}\left(1-i \operatorname{sgn}\langle s, u\rangle \tan \left(\frac{p \pi}{2}\right)\right) \nu(d s) .
$$

Proof. (A) $\Rightarrow$ (B). Let $X_{1}, X_{2}, \ldots$ be i.i.d. with distribution $G$ and $A_{n} \in$ $\mathbb{R}^{d}, B_{n}>0$ such that $\left(S_{n}-A_{n}\right) / B_{n} \rightarrow Z$ weakly where $Z$ is $p$-stable. Let $u \in \mathbb{R}^{d}$. It follows from Feldheim's theorem that $\langle u, Z\rangle$ has a one-dimensional $p$-stable distribution with parameters $\gamma_{u}^{\prime}=\langle u, \mu\rangle, c_{u}^{\prime}=\int_{S^{d-1}}|\langle u, s\rangle|^{p} \nu(d s)$ and

$$
\beta_{u}^{\prime}=\frac{1}{c_{u}^{\prime}} \int_{S^{d-1}}|\langle u, s\rangle|^{p} \operatorname{sgn}(\langle u, s\rangle) \nu(d s) .
$$

The characteristic function $\psi(t u)$ of $\left\langle u, X_{1}\right\rangle$ has a form

$$
\log \psi(t u)=i t \gamma_{u}-|t|^{p} L_{u}\left(\frac{1}{|t|}\right)\left(1-i \beta_{u} \operatorname{sgn}(t) \tan \left(\frac{\pi p}{2}\right)\right)
$$

as in Theorem 1 with some slowly varying function $L_{u}$ and parameters $\gamma_{u}$ and $\beta_{u}$ (we normalize $L_{u}$ so that $c_{u}=1$ ). Hence,

$$
\begin{aligned}
& i t\left(\frac{n \gamma_{u}}{B_{n}}-\frac{\left\langle u, A_{n}\right\rangle}{B_{n}}\right)-|t|^{p} \frac{n}{B_{n}^{p}} L_{u}\left(\frac{B_{n}}{|t|}\right)\left(1-i \beta_{u} \operatorname{sgn}(t) \tan \left(\frac{p \pi}{2}\right)\right) \\
& \quad \rightarrow i t \gamma_{u}^{\prime}-c_{u}^{\prime}|t|^{p}\left(1-i \beta_{u}^{\prime} \operatorname{sgn}(t) \tan \left(\frac{p \pi}{2}\right)\right) .
\end{aligned}
$$

The parameter $\gamma_{u}$ must be linear in $u$ if $p>1$, since $\left(n \gamma_{u}-\left\langle u, A_{n}\right\rangle\right) / B_{n} \rightarrow$ $\langle u, \mu\rangle$ and $n / B_{n} \rightarrow \infty$. In case $p<1, \gamma_{u}$ can be arbitrary since $n / B_{n} \rightarrow 0$. Moreover, $\left(n / B_{n}^{p}\right) L_{u}\left(B_{n}\right)$ converges to $c_{u}^{\prime}$ and $\beta_{u}=\beta_{u}^{\prime}$. Setting $L(t)=$ $\left(1 / c_{u}^{\prime}\right) L_{u}(t)$ for some fixed $u$, we obtain, for $v \in \mathbb{R}^{d}$,

$$
\lim _{n \rightarrow \infty} \frac{L\left(B_{n}\right)}{L_{v}\left(B_{n}\right)}=\lim _{n \rightarrow \infty} \frac{\left(n / B_{n}^{p}\right) L_{u}\left(B_{n}\right)}{c_{u}^{\prime}\left(n / B_{n}^{p}\right) L_{v}\left(B_{n}\right)}=\frac{1}{c_{v}^{\prime}} .
$$

Hence $L_{v}(\lambda) \sim c_{v}^{\prime} L(\lambda)$ as $\lambda \rightarrow \infty$.
$(B) \Rightarrow(A)$. Conversely, if the characteristic function $\psi$ of $G$ is as in (B), then for every $u \in \mathbb{R}^{d}$ the characteristic functions of $Y_{n}^{(u)}=B_{n}^{-1} \sum_{k=1}^{n}\left(\left\langle u, X_{k}\right\rangle-\right.$ $\left\langle A_{n}, u\right\rangle$ ) converge, where $X_{1}, X_{2}, \ldots$ are i.i.d. with distribution $G$, where $B_{n}$ is defined by $n L\left(B_{n}\right)=B_{n}^{p}$ and where $A_{n}=0$ if $p<1$ and $A_{n}=n \mu$ if $p>1$.

It follows that the characteristic functions of $\left(S_{n}-A_{n}\right) / B_{n}$ converge (necessarily to a characteristic function), such that the limit variable $Z$ has all distributions $\langle u, Z\rangle, u \in \mathbb{R}^{d}, p$-stable. Thus, $Z$ is stable itself if $p>1$. In case $p<1$, we note that $Z$ has a characteristic function of the form (1a) with $\mu=0$ and is strictly stable.

If $G$ is a distribution function on $\mathbb{R}^{d}$, we define $G_{u}(\cdot)$ to be the distribution function of $\langle u, Z\rangle$, where $Z$ is a random variable with distribution $G$.

Corollary 2. (A) If a distribution function $G$ on $\mathbb{R}^{d}$ belongs to the domain of attraction of the nondegenerate stablelaw of order 1 , spectral measure $\nu$ and translate $\mu$, then its characteristic function $\psi$ has the form

$$
\begin{align*}
& \operatorname{Relog} \psi(t u)=-t L\left(\frac{1}{t}\right) \int_{S^{d-1}}|\langle u, s\rangle| \nu(d s)+o\left(t L\left(\frac{1}{t}\right)\right) \\
& \operatorname{Im} \log \psi(t u)=t H_{u}\left(\frac{1}{t}\right)+t L\left(\frac{1}{t}\right) \frac{2 C}{\pi} \int_{S^{d-1}}\langle u, s\rangle \nu(d s)+t \gamma_{u}+o\left(t L\left(\frac{1}{t}\right)\right) \tag{7}
\end{align*}
$$

as $t \rightarrow 0^{+} \forall u \in S^{d-1}$, where $L$ is slowly varying at $\infty$,

$$
C=\int_{0}^{\infty}\left(\cos y-\frac{1}{1+y^{2}}\right) \frac{d y}{y}
$$

and

$$
H_{u}(x)=\int_{0}^{x} \frac{v\left(1-G_{u}(v)-G_{u}(-v)\right)}{1+v^{2}} d v
$$

has a representation

$$
\begin{equation*}
H_{u}(\lambda)=\left\langle u, \Gamma_{\lambda}\right\rangle-\frac{2 L(\lambda)}{\pi} \int_{S^{d-1}}\langle u, s\rangle \log (|\langle u, s\rangle|) \nu(d s)-\gamma_{u}+o(L(\lambda)) \tag{8}
\end{equation*}
$$

for some $\Gamma_{\lambda} \in \mathbb{R}^{d}$ and satisfies

$$
\begin{equation*}
H_{u}(k \lambda)-H_{u}(\lambda) \sim \frac{2}{\pi} L(\lambda) \int_{S^{d-1}}\langle u, s\rangle \nu(d s) \log k \tag{9}
\end{equation*}
$$

as $\lambda \rightarrow \infty$.
(B) Let the characteristic function $\psi$ of a distribution $G$ on $\mathbb{R}^{d}$ satisfy (7) for some $\gamma_{u} \in \mathbb{R}$, some finite measure $\nu$ on $S^{d-1}$, some slowly varying function $L$ and some functions $H_{u}$ with representation (8) and satisfying (9). Then $G$ belongs to the domain of attraction of a nondegenerate stable law of order 1.

Proof. (A) As before, let $X_{1}, X_{2}, \ldots$ be i.i.d. with distribution $G$ and $A_{n} \in$ $\mathbb{R}^{d}, B_{n}>0$ such that $\left(S_{n}-A_{n}\right) / B_{n} \rightarrow Z$ weakly, where $Z$ is 1 -stable. Let $u \in \mathbb{R}^{d}$. It follows from Feldheim's theorem that $\langle u, Z\rangle$ has a onedimensional 1-stable distribution with parameters

$$
\begin{aligned}
\gamma_{u}^{\prime} & =\langle u, \mu\rangle-\frac{2}{\pi} \int_{S^{d-1}}\langle u, s\rangle \log (|\langle u, s\rangle|) \nu(d s), \\
c_{u}^{\prime} & =\int_{S^{d-1}}|\langle u, s\rangle| \nu(d s), \quad \beta_{u}^{\prime}=\frac{1}{c_{u}^{\prime}} \int_{S^{d-1}}\langle u, s\rangle \nu(d s) .
\end{aligned}
$$

By Theorem 2, the characteristic function $\psi(t u)$ of $\left\langle u, X_{1}\right\rangle$ has a form

$$
\begin{aligned}
\log \psi(t u)= & -|t| L_{u}\left(\frac{1}{|t|}\right)+i t \gamma_{u}+i t \frac{2 \beta_{u} C}{\pi} L_{u}\left(\frac{1}{|t|}\right) \\
& +i t\left(H_{1 u}\left(\frac{1}{|t|}\right)-H_{2 u}\left(\frac{1}{|t|}\right)\right)+o\left(|t| L_{u}\left(\frac{1}{|t|}\right)\right),
\end{aligned}
$$

where

$$
\begin{aligned}
H_{j u}(\lambda) & =\int_{0}^{\lambda} \frac{x L_{j u}(x)}{1+x^{2}} d x, \\
L_{j u}(x) & = \begin{cases}x\left(1-G_{u}(x)\right), & \text { if } j=1, \\
x G_{u}(-x), & \text { if } j=2,\end{cases}
\end{aligned}
$$

for some parameters $\gamma_{u}, \beta_{u}$ and slowly varying functions $L_{u}$ (normalized so that $c_{u}=1$ ), $L_{j u}$. Also note that, by Theorem 2, $L_{j u}(x)=\left(c_{j u}+o(1)\right) L_{u}(x)$ with $c_{1 u}+c_{2 u}=2 / \pi$. Set $H_{u}=H_{1 u}-H_{2 u}$.

From the assumed convergence of characteristic functions, we have that

$$
\operatorname{Re} n \log \psi\left(\frac{t u}{B_{n}}\right) \sim \frac{n L_{u}\left(B_{n}\right)|t|}{B_{n}} \rightarrow c_{u}^{\prime}|t| .
$$

As in the proof of Corollary 1, there exists a function $L$ so that $c_{v}^{\prime} L \sim L_{v}$ for all $v \in \mathbb{R}^{d}$. Moreover, using (5) $\forall t \in \mathbb{R}$, as $n \rightarrow \infty$,

$$
\begin{aligned}
\operatorname{Im} n & \log \psi\left(\frac{t u}{B_{n}}\right)-\left\langle A_{n}, u\right\rangle \frac{t}{B_{n}} \\
= & \frac{n L_{u}\left(B_{n}\right)}{B_{n}}\left(c_{1 u}-c_{2 u}\right) t \log \frac{1}{|t|} \\
& +t\left(\frac{n \gamma_{u}}{B_{n}}-\frac{\left\langle A_{n}, u\right\rangle}{B_{n}}+\frac{n H_{u}\left(B_{n}\right)}{B_{n}}+\frac{2 C n \beta_{u} L_{u}\left(B_{n}\right)}{\pi B_{n}}\right)+o(1) \\
& \rightarrow t \gamma_{u}^{\prime}+\frac{2 \beta_{u}^{\prime} c_{u}^{\prime} t}{\pi} \log \frac{1}{|t|} .
\end{aligned}
$$

Equating coefficients of $t$, and $t \log 1 /|t|$, we see that

$$
\frac{n L_{u}\left(B_{n}\right)}{B_{n}}\left(c_{1 u}-c_{2 u}\right) \rightarrow \frac{2 \beta_{u}^{\prime} c_{u}^{\prime}}{\pi}
$$

and

$$
\frac{n}{B_{n}}\left(H_{u}\left(B_{n}\right)+\frac{2 C \beta_{u}}{\pi} L_{u}\left(B_{n}\right)+\gamma_{u}-\left\langle u, \frac{A_{n}}{n}\right\rangle\right) \rightarrow \gamma_{u}^{\prime}
$$

as $n \rightarrow \infty$.
Hence, $c_{u}^{\prime}\left(c_{1 u}-c_{2 u}\right)=c_{u}^{\prime} \beta_{u} 2 / \pi=c_{u}^{\prime} 2 \beta_{u}^{\prime} / \pi$ and $\beta_{u}=\beta_{u}^{\prime}$.
To conclude, we determine the conditions for $H_{u}$ and $\gamma_{u}$. Since $c_{u}^{\prime} L \sim L_{u}$ and since $L_{u}$ is slowly varying,

$$
\begin{aligned}
& H_{u}\left(B_{n}\right)+\frac{2 C \beta_{u}^{\prime} c_{u}^{\prime}}{\pi} L\left(B_{n}\right)+\gamma_{u}-\left\langle u, \frac{A_{n}}{n}\right\rangle \\
& \quad-\left\langle u, \frac{B_{n} \mu}{n}\right\rangle+\frac{2 B_{n}}{n \pi} \int_{S^{d-1}}\langle u, s\rangle \log (|\langle u, s\rangle|) \nu(d s)=o\left(\frac{B_{n}}{n}\right)
\end{aligned}
$$

or [because $\beta_{u}^{\prime} c_{u}^{\prime}$ is linear in $u$ and $n L\left(B_{n}\right) \sim B_{n}$ ]

$$
H_{u}\left(B_{n}\right)=\left\langle u, \Gamma_{B_{n}}\right\rangle-\frac{2 L\left(B_{n}\right)}{\pi} \int_{S^{d-1}}\langle u, s\rangle \log (|\langle u, s\rangle|) \nu(d s)-\gamma_{u}+o\left(L\left(B_{n}\right)\right),
$$

where

$$
\Gamma_{B_{n}}=\frac{A_{n}}{n}+\mu L\left(B_{n}\right)-\frac{2 C L\left(B_{n}\right)}{\pi} \int_{S^{d-1}}\langle\cdot, s\rangle \nu(d s) .
$$

We obtain the expansion for $H_{u}(\lambda)$ ( $B_{n} \leq \lambda<B_{n+1}$ ) from

$$
\begin{aligned}
H_{u}(\lambda)-H_{u}\left(B_{n}\right) & =H_{1 u}(\lambda)-H_{1 u}\left(B_{n}\right)-\left[H_{2 u}(\lambda)-H_{2 u}\left(B_{n}\right)\right] \\
& \sim \log \left(\frac{\lambda}{B_{n}}\right)\left(L_{1 u}(\lambda)-L_{2 u}(\lambda)\right)+o(L(\lambda))=o(L(\lambda))
\end{aligned}
$$

and

$$
\begin{aligned}
H_{u}(\lambda) & =H_{u}\left(B_{n}\right)+H_{u}(\lambda)-H_{u}\left(B_{n}\right) \\
& =H_{u}\left(B_{n}\right)+o(L(\lambda)) \\
& =\left\langle u, \Gamma_{B_{n}}\right\rangle-\frac{2 L(\lambda)}{\pi} \int_{S^{d-1}}\langle u, s\rangle \log (|\langle u, s\rangle|) \nu(d s)-\gamma_{u}+o(L(\lambda)),
\end{aligned}
$$

since

$$
1 \leq \frac{\lambda}{B_{n}} \leq \frac{B_{n+1}}{B_{n}} \sim \frac{(n+1) L\left(B_{n+1}\right)}{n L\left(B_{n}\right)} \rightarrow 1 .
$$

Equation (8) follows setting $\Gamma_{\lambda}=\Gamma_{B_{n}}$ if $B_{n} \leq \lambda<B_{n+1}$. Finally, (9) holds because

$$
\begin{aligned}
H_{u}(k \lambda)-H_{u}(\lambda) & \sim \log (k)\left(L_{1 u}(\lambda)-L_{2 u}(\lambda)\right) \\
& \sim \log (k)\left(c_{1 u}-c_{2 u}\right) L_{u}(\lambda) \\
& \sim \log (k)\left(c_{1 u}-c_{2 u}\right) c_{u}^{\prime} L(\lambda) \\
& =\frac{2}{\pi} c_{u}^{\prime} \beta_{u}^{\prime} \log (k) L(\lambda) .
\end{aligned}
$$

(B) Conversely, if the characteristic function $\psi$ of $G$ is as in (B), then for every $u \in \mathbb{R}^{d}$ the characteristic functions of

$$
Y_{n}^{(u)}=B_{n}^{-1} \sum_{k=1}^{n}\left(\left\langle u, X_{k}\right\rangle-\left\langle A_{n}, u\right\rangle\right)
$$

converge, where $X_{1}, X_{2}, \ldots$ are i.i.d. with distribution $G$, where $B_{n}$ is defined by $n L\left(B_{n}\right)=B_{n}$ and where

$$
A_{n}=n \Gamma_{B_{n}}+\frac{2 C n L\left(B_{n}\right)}{\pi} \int_{S^{d-1}}\langle\cdot, s\rangle \nu(d s) .
$$

Let $c_{u}^{\prime}=\int_{S^{d-1}}|\langle u, s\rangle| \nu(d s)$ be defined as before. We have that

$$
\begin{aligned}
& \log \left(\psi\left(\frac{t u}{B_{n}}\right)^{n} \exp \left[-\frac{i t\left\langle u, A_{n}\right\rangle}{B_{n}}\right]\right) \\
& \quad \rightarrow-|t| c_{u}^{\prime}-i t \frac{2}{\pi} \int_{S^{d-1}}\langle u, s\rangle \log |\langle t u, s\rangle| \nu(d s)
\end{aligned}
$$

Example. Let $0<p<2, \nu \in \mathscr{P}\left(S^{d-1}\right)$ be nondegenerate, and let $L$ be slowly varying at $\infty$.

If $Y \in \operatorname{DA}(p, 1), Y>0$ with tails given by $P(Y>\lambda)=2 L(\lambda) / \pi \lambda^{p}$ and $Z$ is a $\nu$-distributed random variable on $S^{d-1}$ independent of $Y$, then $X:=Y Z$ is in the domain of attraction of a nondegenerate stable law of order $p$ on $\mathbb{R}^{d}$ and with spectral measure $\nu$.

This follows from (and illustrates) Corollaries 1 and 2. Indeed, using the notation $\psi_{U}(u):=-\log \left(E[\exp (i\langle U, u\rangle)]\right.$, we have that, for $u \in S^{d-1}$ and $t>0$,

$$
\begin{aligned}
\psi_{X}(t u) & =E\left(\psi_{Y}(\langle Z, t u\rangle)+O\left(\psi_{Y}(\langle Z, t u\rangle)^{2}\right)\right) \\
& =E\left(\psi_{Y}(\langle Z, t u\rangle)\right)+o\left(t^{p} L(1 / t)\right)
\end{aligned}
$$

as $t \rightarrow 0$, whence, by Ibragimov and Linnik (1971) for $p \neq 1$,

$$
\begin{aligned}
\psi_{X}(t u)= & i t \gamma\langle u, E(Z)\rangle \\
& -t^{p} L\left(\frac{1}{t}\right) \int_{S^{d-1}}|\langle u, s\rangle|^{p}\left(1-i \operatorname{sgn}(\langle s, u\rangle) \tan \left(\frac{p \pi}{2}\right)\right) \nu(d s) \\
& +o\left(t^{p} L\left(\frac{1}{t}\right)\right)
\end{aligned}
$$

as $t \rightarrow 0$, and, by Theorem 2 for $p=1$,

$$
\begin{aligned}
\operatorname{Re} \psi_{X}(t u)= & -t L\left(\frac{1}{t}\right) \int_{S^{d-1}}|\langle s, u\rangle| d \nu(s)+o\left(t L\left(\frac{1}{t}\right)\right) \\
\operatorname{Im} \psi_{X}(t u)= & t \gamma\langle u, E(Z)\rangle+t\left(H\left(\frac{1}{t}\right)+\frac{2 C}{\pi} L\left(\frac{1}{t}\right)\right) \int_{S^{d-1}}\langle s, u\rangle d \nu(s) \\
& +t L\left(\frac{1}{t}\right) \frac{2}{\pi} \int_{S^{d-1}}\langle s, u\rangle \log \frac{1}{|\langle s, u\rangle|} d \nu(s)+o\left(t L\left(\frac{1}{t}\right)\right)
\end{aligned}
$$

as $t \rightarrow 0$, where

$$
H(\lambda):=\int_{0}^{\lambda} \frac{2 x L(x) d x}{\pi\left(1+x^{2}\right)}
$$

and where

$$
\gamma:=E\left(\frac{Y}{1+Y^{2}}+\int_{0}^{Y} \frac{2 u^{2}}{\left(1+u^{2}\right)^{2}} d u\right)
$$

If, in the example, $Y$ was not chosen positive, but satisfying (2) with constants $c, c_{1}, c_{2}$, then the spectral measure of $X$ is given by

$$
\nu^{*}(A)=c_{1} \nu(A)+c_{2} \nu(-A), \quad A \in \mathscr{B}\left(S^{d-1}\right)
$$

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