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On uniqueness of moving average representations of heavy-tailed stationary processes

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SUMMARY

We prove the uniqueness of linear i.i.d. representations of heavy-tailed processes whose distribution belongs to the domain of attraction of an α -stable law, with $\alpha < 2$. This shows the possibility to identify nonparametrically both the sequence of two-sided moving average coefficients and the distribution of the heavy-tailed i.i.d. process.

Some key words: α -stable distribution; Domain of attraction; Infinite moving average; Linear process; Mixed causal/noncausal process; Nonparametric identification; Unobserved component model.

1. INTRODUCTION

By definition, a real linear process (X_t) can be written as a two-sided moving average of a strong white noise:

$$X_t = \sum_{j=-\infty}^{\infty} a_j \epsilon_{t-j}, \quad t \in \mathbb{Z}, \quad (1)$$

where (ϵ_t) is a sequence of independent and identically distributed (i.i.d.) real random variables, and (a_j) is a sequence of Moving-Average (MA) coefficients.

Provided that such a process exists, an important issue is the uniqueness of the MA representation in (1). Uniqueness of the linear representation of non-Gaussian i.i.d. processes plays an important role in time series, for instance in the analysis of time reversibility (see Hallin, Lefèvre and Puri (1988), Breidt and Davis (1992)). From a statistical point of view, this problem is equivalent to the nonparametric identification of the sequence (a_i) and of the common distribution of the variables ϵ_t . See Rosenblatt (2000) for a review of statistical applications of non-Gaussian linear processes.

While Gaussian linear processes generally admit several moving average representations, non-Gaussian linear processes have been shown to admit an essentially unique MA representation, under different regularity conditions on the MA coefficients and the moments of the i.i.d. process. The literature dealing with uniqueness of MA representations of non-Gaussian processes includes Lii and Rosenblatt (1982), Findley (1986, 1990), Cheng (1992), Breidt and Davis (1992). The sufficient uniqueness conditions obtained in these articles are summarized in Table 1. They require at least finite variance of the independent process (ϵ_t) , except in Breidt and Davis (1992) (see their Remark 2, p. 286). In the present article, we provide a direct proof for the identifi-

cation of the MA representation when the errors belong to the domain of attraction of a stable law, with index $\alpha < 2$. This case is especially important in view of the increased interest, in the recent time series literature, in heavy-tailed distributions. In particular, it has been noted that mixed causal/noncausal processes of the form (1) were adequate for modeling the speculative bubbles observed on fat-tailed financial series, for instance when the noise and the observed process do not have a finite expectation (see Gouriéroux and Zakoïan (2013)). Finally, we extend the identification result to the multivariate framework, especially to unobserved component models.

2. UNIQUENESS RESULTS FOR UNIVARIATE PROCESSES

We start by recalling conditions ensuring the existence of a MA representation of a strictly stationary process.

2.1. Existence

The existence of infinite MA processes, as in (1), is generally established under the assumption that the variables ϵ_t belong to L^1 or L^2 . An extension was established by Cline (Theorem 2.1, 1983), who gave conditions for the existence of infinite MA with i.i.d. random variables admitting regularly varying tails. The following is a straightforward consequence of Proposition 13.3.1 in Brockwell and Davis (1991).

PROPOSITION 1. *Consider the two-sided moving average*

$$X_t = \sum_{j=-\infty}^{\infty} \alpha_j Z_{t-j}, \quad (2)$$

written on a process (Z_t) such that, for some $s \in (0, 1)$:

$$\sup_t E|Z_t|^s < \infty \quad \text{and} \quad \sum_{h=-\infty}^{\infty} |\alpha_h|^s < \infty.$$

Then, the series X_t converges absolutely with probability one. If, in addition, (Z_t) is a strictly stationary process, then (X_t) is also strictly stationary.

We now turn to the uniqueness problem, first in the case where $X_t = \epsilon_t^*$ is an i.i.d. process.

2.2. Uniqueness of the MA representation of an heavy-tailed i.i.d. process

Let (ϵ_t) and (ϵ_t^*) denote two i.i.d. processes. Suppose that

$$\epsilon_t^* = \sum_{j=-\infty}^{\infty} a_j \epsilon_{t-j}, \quad t \in \mathbb{Z}, \quad (3)$$

and, for some $s \in (0, 1)$,

$$E|\epsilon_t|^s < \infty \quad \text{and} \quad \sum_{j=-\infty}^{\infty} |a_j|^s < \infty. \quad (4)$$

By Proposition 1, the latter conditions entail the existence of the infinite MA appearing in the right-hand side of (3). Note also that this entails $E|\epsilon_t^*|^s < \infty$, using the inequality $(x + y)^s \leq x^s + y^s$ for $x, y \geq 0$. Let Ψ, Ψ^* denote the characteristic functions of ϵ and ϵ^* , respectively.

The characteristic function of $(\epsilon_t^*, \epsilon_{t-h}^*)$ with $h \neq 0$ is, for $u, v \in \mathbb{R}$,

$$E\left(e^{i(u\epsilon_t^* + v\epsilon_{t-h}^*)}\right) = \Psi^*(u)\Psi^*(v) = E\left(\prod_{j=-\infty}^{\infty} e^{i(ua_j + va_{j+h})\epsilon_{t-j}}\right) = \prod_{j=-\infty}^{\infty} \Psi(ua_j + va_{j+h}).$$

The last equality holds by the dominated convergence theorem, which applies because the products $\prod_{j=-n}^n e^{i(ua_j + va_{j+h})\epsilon_{t-j}}$ have unit norm.

Hence, for $u, v \in \mathbb{R}$,

$$\log |\Psi^*(u)| + \log |\Psi^*(v)| = \sum_{j=-\infty}^{\infty} \log |\Psi(ua_j + va_{j+h})|. \quad (5)$$

For expository purpose, we first consider stable variables.

a) Stable variables

Let us recall that the distribution of a stable variable X with parameters $(\alpha, \beta, \sigma, \mu)$, denoted $X \sim \mathcal{S}(\alpha, \beta, \sigma, \mu)$, has the characteristic function

$$\Psi(s) = \begin{cases} \exp[-\sigma^\alpha |s|^\alpha \{1 - i\beta (\text{sign } s) \tan(\frac{\pi\alpha}{2})\} + i\mu s], & \text{if } \alpha \neq 1, \\ \exp[-\sigma |s| \{1 + i\beta (\text{sign } s) \frac{2}{\pi} \log |s|\} + i\mu s], & \text{if } \alpha = 1. \end{cases}$$

For a stable distribution with index of stability $\alpha \in (0, 2)$, the moment condition in (4) is satisfied for $s \in (0, \min\{\alpha, 1\})$. See Samorodnitsky and Taqqu (1994) for further details on stable variables.

PROPOSITION 2. *If (ϵ_t) and (ϵ_t^*) have stable laws, $\epsilon_t \sim \mathcal{S}(\alpha, \beta, \sigma, \mu)$ and $\epsilon_t^* \sim \mathcal{S}(\alpha_*, \beta_*, \sigma_*, \mu_*)$ with $\alpha, \alpha_* < 2$, if $\sum_{j=-\infty}^{\infty} |a_j|^s < \infty$ for $s \in (0, \min\{\alpha, 1\})$, then all the MA coefficients a_j 's except one are equal to zero in the MA representation (3).*

b) Variables in the domain of attraction of a stable law

By definition, a variable X belongs to the domain of attraction of an α -stable law, denoted $X \in \text{DA}(\alpha)$, with $\alpha < 2$, if for any i.i.d. sequence (X_i) , with $X_1 \stackrel{d}{=} X$, there exist sequences of constants $A_n \in \mathbb{R}$ and $B_n > 0$ such that

$$\frac{X_1 + \cdots + X_n - A_n}{B_n} \xrightarrow{d} Z, \quad (6)$$

where \xrightarrow{d} denotes the convergence in distribution, and the distribution of Z is α -stable. The tails of variables belonging to $\text{DA}(\alpha)$ have the following characterization: $X \in \text{DA}(\alpha)$ if and only if its cumulative distribution function is such that:

$$F(-x) \sim \frac{c_1 L(x)}{x^\alpha}, \quad 1 - F(x) \sim \frac{c_2 L(x)}{x^\alpha}, \quad \text{when } x \rightarrow \infty,$$

where L is a slowly varying function¹ and $c_1, c_2 \geq 0$, $c_1 + c_2 \neq 0$ (see for instance Embrechts, Klüppelberg, Mikosch (1997), Theorem 2.2.8). Hence, if $X \in \text{DA}(\alpha)$, $E|X|^s < \infty$ for $s \in (0, \alpha)$. Note that we can have $E|X|^\alpha < \infty$ which is not possible for α -stable variables.

Let us start with a simple extension of Proposition 2.

¹ that is, L is measurable positive, defined on $(0, \infty)$, such that $\lim_{x \rightarrow \infty} L(tx)/L(x) = 1$ for all $t > 0$.

PROPOSITION 3. If $\epsilon_t^* \in DA(\alpha_*)$ and ϵ_t has a stable distribution with index α , with $\alpha, \alpha_* < 2$, if $\sum_{j=-\infty}^{\infty} |a_j|^s < \infty$ for $s < \alpha$, then all the MA coefficients a_j 's except one are equal to zero in (3).

Let us now extend the identification result, when both i.i.d. processes (ϵ_t) and (ϵ_t^*) are in the attraction domain of a stable law.

PROPOSITION 4. If $\epsilon_t \in DA(\alpha)$ and $\epsilon_t^* \in DA(\alpha_*)$ are such that

$$|E(e^{iu\epsilon_t})| = e^{-|u|^\alpha L_1(1/u)}, \quad |E(e^{iu\epsilon_t^*})| = e^{-|u|^{\alpha_*} L_1^*(1/u)}, \quad (7)$$

with $\alpha, \alpha_* < 2$, where L_1, L_1^* are slowly varying functions, if for any $\tau > 0$ there exist constants $M, K > 0$ such that for all $t > 0$,

$$\sup_{|x| > M} \left| \frac{L_1(tx)}{L_1(x)} \right| \leq K(1 + t^\tau), \quad (8)$$

and if $\sum_{j=-\infty}^{\infty} |a_j|^s < \infty$ for $s < \alpha$, then all the MA coefficients a_j 's except one are equal to zero in (3).

2.3. Uniqueness of the MA representation of an heavy-tailed linear process

We now study the uniqueness of representation (1). Assume that

$$X_t = \sum_{j=-\infty}^{\infty} a_j \epsilon_{t-j} = \sum_{j=-\infty}^{\infty} a_j^* \epsilon_{t-j}^*, \quad t \in \mathbb{Z} \quad (9)$$

where (ϵ_t) and (ϵ_t^*) are two sequences of i.i.d. real random variables, and $(a_j), (a_j^*)$ are two sequences of MA coefficients with, for some $s \in (0, 1)$:

$$\sum_{h=-\infty}^{\infty} |a_h|^s < \infty \quad \text{and} \quad \sum_{h=-\infty}^{\infty} |a_h^*|^s < \infty. \quad (10)$$

The random series in (9) converging a.s., we may define $a(B)\epsilon_t := \sum_{j=-\infty}^{\infty} a_j \epsilon_{t-j}$ and $a^*(B)\epsilon_t^* := \sum_{j=-\infty}^{\infty} a_j^* \epsilon_{t-j}^*$ where B stands for the backward shift operator. The next proposition shows that under appropriate assumptions on the i.i.d. sequence and the MA coefficients, the representation (1) of X_t is essentially unique.

PROPOSITION 5. Assume that $\epsilon_t \in DA(\alpha)$ and $\epsilon_t^* \in DA(\alpha_*)$, satisfying conditions (7)-(8), that $a^*(B)$ is invertible, with $\{a^*(B)\}^{-1} = \sum_{j=-\infty}^{\infty} c_j B^j$ such that $\sum_{j=-\infty}^{\infty} |c_j|^s < \infty$, and that $\sum_{j=-\infty}^{\infty} |a_j|^s < \infty$, for $s < \alpha$. Then, if (9)-(10) hold, we have

$$\epsilon_t^* = c\epsilon_{t-\ell} \quad \text{and} \quad a_j = ca_{j+\ell}^*, \quad \forall j \in \mathbb{Z}$$

for some constants $c \in \mathbb{R}$ and $\ell \in \mathbb{Z}$.

From Proposition 5, we see that the MA representation is identifiable up to a change of scale and a drift of the time index on ϵ_t .

3. UNIQUENESS RESULTS FOR MULTIVARIATE HEAVY-TAILED LINEAR PROCESSES

The literature on the identification of multivariate MA processes is rather limited. Sets of conditions for identification have been derived in Chan, Ho and Tong (2006), when the p -dimensional observed process and the q -dimensional noise processes are square integrable, when

$p \geq q$ (i.e. an order condition), and when the q components of the noise are i.i.d. (see their Condition 3). The order condition is required in the proof to be able to apply the approach by Cheng (1992, 1999).

This section will complete such results by considering processes with fat tails. We start by an extension of Proposition 3.

3.1. Uniqueness of the decomposition of an heavy-tailed i.i.d. process as a sum of independent stable MA

The next result extends Proposition 3 by considering the decomposition of a non-Gaussian stable one-dimensional noise as the sum of independent stable MA processes:

$$\epsilon_t^* = \sum_{k=1}^K \sum_{j=-\infty}^{\infty} a_{k,j} \epsilon_{k,t-j}, \quad t \in \mathbb{Z} \quad (11)$$

where $(\epsilon_{k,t})$ are K independent sequences of stable i.i.d. real random variables, and the $\{a_{k,j}; j \in \mathbb{Z}\}$ are K sequences of real numbers such that, for some $s \in (0, 1)$,

$$0 < \sum_{j=-\infty}^{\infty} |a_{k,j}|^s < \infty.$$

PROPOSITION 6. *If $\epsilon_t^* \in DA(\alpha_*)$ with $\alpha_* < 2$, the variables $\epsilon_{k,t}$ have stable laws, $\epsilon_{k,t} \sim \mathcal{S}(\alpha_k, \beta_k, \sigma_k, \mu_k)$ with $\alpha_k \in (s, 2)$ for $k = 1, \dots, K$, then in each sequence $\{a_{k,j}; j \in \mathbb{Z}\}$ all the $a_{k,j}$'s except one are equal to zero in (11).*

3.2. Uniqueness of the MA representation of a vector process with fat tails

We now show that, under appropriate conditions, a vector process (\mathbf{X}_t) has at most one representation of the form

$$\mathbf{X}_t = \mathbf{A}(B)\epsilon_t, \quad (12)$$

where (ϵ_t) is a vector i.i.d. process, with independent components following stable distributions with distinct tail indices. More precisely, consider a p -dimensional process $\mathbf{X}_t = (X_{1t}, \dots, X_{pt})'$ admitting two MA(∞) representations given by

$$\mathbf{X}_t = \mathbf{A}(B)\epsilon_t = \mathbf{A}^*(B)\epsilon_t^*, \quad t \in \mathbb{Z}, \quad (13)$$

where $(\epsilon_t), (\epsilon_t^*)$ are two K -dimensional i.i.d. processes with stable and independent components, and $\mathbf{A}(B) = (a_{i,k}(B))_{i=1, \dots, p; k=1, \dots, K}$, $\mathbf{A}^*(B) = (a_{i,k}^*(B))_{i=1, \dots, p; k=1, \dots, K}$ are infinite-order lag polynomials with matrix coefficients. Let $a_{i,k}(B) = \sum_{j=-\infty}^{\infty} a_{i,k,j} B^j$ and $a_{i,k}^*(B) = \sum_{j=-\infty}^{\infty} a_{i,k,j}^* B^j$. Assume that, for some $s \in (0, 1)$,

$$0 < \sum_{j=-\infty}^{\infty} |a_{i,k,j}|^s < \infty, \quad 0 < \sum_{j=-\infty}^{\infty} |a_{i,k,j}^*|^s < \infty.$$

PROPOSITION 7. *Assume that (ϵ_t) and (ϵ_t^*) are two K -dimensional i.i.d. processes satisfying (13), with stable and independent components $\epsilon_{k,t}$ and $\epsilon_{k,t}^*$ having stable laws, $\epsilon_{k,t} \sim \mathcal{S}(\alpha_k, \beta_k, \sigma_k, \mu_k)$ and $\epsilon_{k,t}^* \sim \mathcal{S}(\alpha_k^*, \beta_k^*, \sigma_k^*, \mu_k^*)$, with $\alpha_k, \alpha_k^* \in (s, 2)$ for $k = 1, \dots, K$ and $\alpha_1 < \dots < \alpha_K$. Suppose that, for any i, k , the lag polynomial $a_{i,k}^*(B)$ is invertible, with $\{a_{i,k}^*(B)\}^{-1} = \sum_{j=-\infty}^{\infty} c_{i,k,j} B^j$ such that $\sum_{j=-\infty}^{\infty} |c_{i,k,j}|^s < \infty$. Then, for $i = 1, \dots, p$, and*

$k = 1, \dots, K$

$$a_{i,k,j} = c_{i,k} a_{i,k,j+\ell_{i,k}}^*, \quad \forall j \in \mathbb{Z},$$

for some constants $c_{i,k} \in \mathbb{R}$ and $\ell_{i,k} \in \mathbb{Z}$.

4. CONCLUDING REMARKS

The identification results of this paper can be used in both parametric and nonparametric analyses of linear processes. The proofs also suggest a nonparametric estimation method, for a linear process X_t based on an i.i.d. sequence of heavy-tailed variables ϵ_t , as in (1), by the following steps.

i) Estimation of α . Since $X_t \in \text{DA}(\alpha)$ when $\epsilon_t \in \text{DA}(\alpha)$ (see for instance Embrechts, Klüppelberg, Mikosch (1997), Theorem A3.26), the tail index α can be consistently estimated by a standard approach. For instance, the Hill estimator can be used (for its main properties under various assumptions see Embrechts, Klüppelberg, Mikosch (1997) Theorem 6.4.6).

ii) Estimation of the MA coefficients. The analysis of the joint characteristic function of (X_t, X_{t-h}) for $h \neq 0$ and $v = tu$ in a neighborhood of $u = 0$, provides estimation of $\sum_{j=-\infty}^{\infty} |a_j + ta_{j+h}|^\alpha$ when $\epsilon_t \sim \mathcal{S}(\alpha, \beta, \sigma, \mu)$. Estimates of (a finite number of) MA coefficients can be deduced by truncating the infinite sum and solving the moment restrictions.

iii) Estimation of the distribution of ϵ_t . Having estimated the lag polynomial $a(B)$, approximations of the variables ϵ_t can be deduced by computing the residuals $\hat{\epsilon}_t = \hat{a}(B)^{-1} X_t$, where the X_u for $u < 0$ are replaced by 0. Then, the common distribution of the ϵ_t 's will be estimated from the residuals.

APPENDIX 1

Uniqueness of non-Gaussian MA representations in the literature

Table 1 gives a summary of the main conditions obtained in the literature for the uniqueness of the MA representation in (1). Apart from such conditions, it is generally also assumed that the spectral density of X_t is positive almost everywhere.

Note that the condition given by Breidt and Davis (Remark 2, p.386, 1992), as an alternative proof to their Proposition 3.1, is based on the following result:

THEOREM 1 (KAGAN ET AL. (1973, P.94)). *Let ϵ_j be a sequence of independent random variables, and $\{a_j\}, \{b_j\}$ be two sequences of real constants such that*

- i) *The sequences $\{a_j/b_j : a_j b_j \neq 0\}$ and $\{b_j/a_j : a_j b_j \neq 0\}$ are both bounded; $\sum a_j \epsilon_j$ and $\sum b_j \epsilon_j$ converge a.s. to random variables U and V respectively;*
- ii) *U and V are independent.*

Then, for every j such that $a_j b_j \neq 0$, ϵ_j is normally distributed.

In the context of MA processes, this result can be used as follows. Suppose that (3)-(4) hold, where (ϵ_t) and (ϵ_t^*) are two i.i.d. processes. Then the variables

$$\epsilon_t^* = \sum_{j=-\infty}^{\infty} a_j \epsilon_{t-j} \quad \text{and} \quad \epsilon_{t-h}^* = \sum_{j=-\infty}^{\infty} a_{j+h} \epsilon_{t-j}$$

are independent. It follows that if the sequences

$$\{a_j/a_{j+h} : a_j a_{j+h} \neq 0\} \quad \text{and} \quad \{a_{j+h}/a_j : a_j a_{j+h} \neq 0\} \quad \text{are both bounded,} \quad (\text{A1})$$

Table 1. *Uniqueness conditions for the non-Gaussian MA representation (1)*

Papers	Condition on ϵ_t	Condition on MA coefficients
Lii and Rosenblatt (1982)	All moments finite	$\sum j a_j < \infty$
Findley (1986)	All moments finite	$\sum a_j^2 < \infty$
Findley (1990)	Finite non-zero r th cumulant	$\sum a_j^2 < \infty$
Cheng (1992)	Finite variance	$\sum a_j^2 < \infty$
Breidt and Davis (Prop. 3.1, 1992)	Finite variance	Fractionally Integrated ARMA
Breidt and Davis (Rem. (2) p. 386, 1992)	$E \log^+ \epsilon_t < \infty$	Fractionally Integrated ARMA

then either the variables ϵ_t and ϵ_t^* are normally distributed, or all the a_j 's except one are equal to zero. The conclusion thus coincides with that of Proposition 4, but is obtained under (a) less restrictive distributional assumptions on the i.i.d. processes, but (b) more restrictive assumptions on the sequence of coefficients a_j . Breidt and Davis (Remark 2, p. 386, 1992) showed that (A1) is satisfied for a MA polynomial of the form $\mathcal{A}(z) = \Psi(z)/\Phi(z) = \sum_{j=0}^{\infty} a_j z^j$, where $\Psi(z)$ and $\Phi(z)$ are finite-order polynomials under standard assumptions. For more general processes, the boundedness assumption can be very restrictive. For instance, it precludes sequences recursively defined by $a_j = 0$ for $j \leq 0$ and $a_j = \lambda_j a_{j-1}$ for $j \geq 1$, where (λ_j) is a sequence converging to zero or to infinity.

APPENDIX 2

Proof of Proposition 2

In view of (5) with $v = 0$, we have, for all $u \in \mathbb{R}$,

$$-\sigma_*^{\alpha_*} |u|^{\alpha_*} = - \sum_{j=-\infty}^{\infty} \sigma^\alpha |a_j|^\alpha |u|^\alpha.$$

Thus $\alpha = \alpha_*$ and, without loss of generality, we can take $\sigma = \sigma_*$ and $\sum_{j=-\infty}^{\infty} |a_j|^\alpha = 1$. We then have, by taking $v = tu$ in (5),

$$\sum_{j=-\infty}^{\infty} |a_j + ta_{j+h}|^\alpha = 1 + |t|^\alpha, \quad \forall t \in \mathbb{R}. \quad (\text{A2})$$

The right-hand side of this equality is a differentiable function of t except at $t = 0$. Let us study the differentiability of the left-hand side. Suppose that there exists j_0 such that $a_{j_0} \neq 0$ and $a_{j_0+h} \neq 0$. The function $t \mapsto |a_{j_0} + ta_{j_0+h}|^\alpha$ is everywhere differentiable except at $t_0 = -a_{j_0}/a_{j_0+h}$, with $t_0 \neq 0$. The contribution of all terms of the infinite sum displaying non differentiability at t_0 is

$$\sum_{j \in J} |a_j + ta_{j+h}|^\alpha = |t - t_0|^\alpha \sum_{j \in J} |a_{j+h}|^\alpha,$$

where $J = \{j \in \mathbb{Z} \mid a_j + t_0 a_{j+h} = 0\}$. This contribution must vanish; otherwise, we would have in the right-hand side of (A2) a function which is differentiable at t_0 while the left-hand side is not. Hence $\sum_{j \in J} |a_{j+h}|^\alpha = 0$, in contradiction with the fact that this sum is bounded below by $|a_{j_0+h}|^\alpha \neq 0$. Therefore, there exists no integer j_0 such that $a_{j_0} \neq 0$ and $a_{j_0+h} \neq 0$, which establishes the property. \square

Proof of Proposition 3

By Ibragimov and Linnik (Theorem 2.6.5, 1971), for a variable $\epsilon^* \in \text{DA}(\alpha_*)$, where the limiting stable distribution has location parameter $\mu = 0$, we have

$$|E(e^{iu\epsilon^*})| = e^{-|u|^{\alpha_*} L_1^*(1/u)}, \quad (\text{A3})$$

in the neighborhood of the origin, for some slowly varying function L_1^* . This equality also holds if $\mu \neq 0$, the limit law in (6) being uniquely determined up to positive affine transformations. By (5) with $v = 0$, we thus have

$$-|u|^{\alpha_*} L_1^* \left(\frac{1}{|u|} \right) = - \sum_{j=-\infty}^{\infty} \sigma^\alpha |a_j|^\alpha |u|^\alpha, \quad (\text{A4})$$

in the neighborhood of 0. Thus $\alpha = \alpha_*$ and L_1^* is constant in the neighborhood of 0. It is not restrictive to assume $\sigma^\alpha = L_1^*(1/|u|)$ in the neighborhood of 0 and $\sum_{j=-\infty}^{\infty} |a_j|^\alpha = 1$. We thus are lead to the proof of Proposition 2. □

Proof of Proposition 4

By (5) with $v = 0$ and (A3), we have in the neighborhood of 0,

$$-|u|^{\alpha_*} L_1^* \left(\frac{1}{|u|} \right) = - \sum_{\substack{j=-\infty \\ a_j \neq 0}}^{\infty} |a_j|^\alpha |u|^\alpha L_1 \left(\frac{1}{|a_j u|} \right). \quad (\text{A5})$$

i) Let us first show that $\alpha_* = \alpha$. If $\alpha_* < \alpha$, the left-hand side term in

$$|u|^{\alpha_* - \alpha} L_1^* \left(\frac{1}{|u|} \right) = - \sum_{\substack{j=-\infty \\ a_j \neq 0}}^{\infty} |a_j|^\alpha L_1 \left(\frac{1}{|a_j u|} \right),$$

tends to 0 when $u \rightarrow 0$, while the right-hand side term does not. We also get a contradiction if $\alpha_* < \alpha$, hence $\alpha = \alpha_*$.

ii) Now we will show that

$$L_1^* \left(\frac{1}{|u|} \right) \sim \left(\sum_{j=-\infty}^{\infty} |a_j|^\alpha \right) L_1 \left(\frac{1}{|u|} \right), \quad \text{when } u \rightarrow 0. \quad (\text{A6})$$

We will use the next property (see for instance Embrechts, Klüppelberg, Mikosch (1997), Theorem A3.2).

LEMMA 1. *If L is a slowly varying function, for any $a, b > 0$, the convergence $L(tx)/L(x) \rightarrow 1$ when $x \rightarrow \infty$ is uniform on the segment $[a, b]$.*

Using Lemma 1 and (A5) with $\alpha = \alpha^*$ we have, in the neighborhood of 0, for some constants $a, b > 0$ to be chosen later,

$$\begin{aligned} & \left| \frac{L_1^* \left(\frac{1}{|u|} \right)}{L_1 \left(\frac{1}{|u|} \right)} - \sum_{j=-\infty}^{\infty} |a_j|^\alpha \right| \\ & \leq \sum_{\substack{j=-\infty \\ a_j \neq 0}}^{\infty} |a_j|^\alpha \left| \frac{L_1 \left(\frac{1}{|a_j u|} \right)}{L_1 \left(\frac{1}{|u|} \right)} - 1 \right| \\ & \leq \left(\sum_{j=-n}^n |a_j|^\alpha \right) \sup_{t \in [a, b]} \left| \frac{L_1 \left(\frac{t}{|u|} \right)}{L_1 \left(\frac{1}{|u|} \right)} - 1 \right| + \sum_{\substack{|j| > n \\ a_j \neq 0}} |a_j|^\alpha \left| \frac{L_1 \left(\frac{1}{|a_j u|} \right)}{L_1 \left(\frac{1}{|u|} \right)} \right| + \sum_{|j| > n} |a_j|^\alpha \\ & := S_{1,n}(u) + S_{2,n}(u) + S_{3,n}. \end{aligned}$$

We have, in the neighborhood of 0, for $\tau = \alpha - s$,

$$S_{2,n}(u) = \sum_{\substack{|j|>n \\ a_j \neq 0}} |a_j|^s |a_j|^\tau \left| \frac{L_1\left(\frac{1}{|a_j u|}\right)}{L_1\left(\frac{1}{|u|}\right)} \right| \leq \sum_{\substack{|j|>n \\ a_j \neq 0}} |a_j|^s |a_j|^\tau K (1 + |a_j|^{-\tau}),$$

which is smaller than an arbitrarily small $\varsigma > 0$ for n sufficiently large, $n > N$ say. The same property holds for $S_{3,n}$. Now let $[a, b] = [\min_{|j| \leq N} \{|a_j|^{-1}, a_j \neq 0\}, \max_{|j| \leq N} \{|a_j|^{-1}, a_j \neq 0\}]$ in $S_{1,n}(u)$. We have

$$S_{1,n}(u) \leq \left(\sum_{j=-\infty}^{\infty} |a_j|^\alpha \right) \sup_{t \in [a,b]} \left| \frac{L_1\left(\frac{t}{|u|}\right)}{L_1\left(\frac{1}{|u|}\right)} - 1 \right|,$$

which tends to zero using the uniform convergence of L_1 .

Thus we have established (A6) and, without generality loss, we can make a scale change on the a_j 's to get, for $u \neq 0$,

$$L_1^*\left(\frac{1}{|u|}\right) = \left(\sum_{j=-\infty}^{\infty} |a_j|^\alpha \right) L_1\left(\frac{1}{|u|}\right). \quad (\text{A7})$$

iii) Using (A7), Equation (5) can now be written, for u, v in the neighborhood of 0, $uv \neq 0$, as

$$\begin{aligned} |u|^\alpha L_1^*\left(\frac{1}{|u|}\right) + |v|^\alpha L_1^*\left(\frac{1}{|v|}\right) &= \sum_{\substack{j=-\infty \\ a_j u + a_{j+h} v \neq 0}}^{\infty} |a_j u + a_{j+h} v|^\alpha L_1\left(\frac{1}{|a_j u + a_{j+h} v|}\right) \\ &= \sum_{j=-\infty}^{\infty} |a_j u|^\alpha L_1\left(\frac{1}{|u|}\right) + |a_{j+h} v|^\alpha L_1\left(\frac{1}{|v|}\right). \end{aligned}$$

Hence, for $v = tu$, $t > 0$, we get, for u in the neighborhood of 0,

$$\begin{aligned} &\sum_{\substack{j=-\infty \\ a_j + ta_{j+h} \neq 0}}^{\infty} |a_j + ta_{j+h}|^\alpha L_1\left(\frac{1}{|a_j + ta_{j+h}| |u|}\right) \\ &= \sum_{j=-\infty}^{\infty} |a_j|^\alpha L_1\left(\frac{1}{|u|}\right) + |ta_{j+h}|^\alpha L_1\left(\frac{1}{t|u|}\right). \end{aligned}$$

By the same argument as before, using the uniform convergence of the slowly varying function L_1 (see Lemma 1), the asymptotic behavior of these sums in the neighborhood of $u = 0$ is the same if we replace terms of the form $L_1\left(\frac{1}{x|u|}\right)$, with $x > 0$, by $L_1\left(\frac{1}{|u|}\right)$. We thus have, for any $t > 0$,

$$\sum_{j=-\infty}^{\infty} |a_j + ta_{j+h}|^\alpha = \sum_{j=-\infty}^{\infty} (|a_j|^\alpha + |ta_{j+h}|^\alpha),$$

and we are led to the discussion of Proposition 2. □

Proof of Proposition 5

We use Theorem 2.1 in Kokoszka (1996), which establishes sufficient conditions for two infinite-order lag polynomials whose coefficients may not be absolutely summable to commute. From the assumptions made on the coefficients of the polynomials $a(B)$ and $\{a^*(B)\}^{-1}$, we have $\epsilon_t^* = \{a^*(B)\}^{-1} a(B) \epsilon_t$, where the random series $\{a^*(B)\}^{-1} a(B) \epsilon_t$ converges absolutely a.s. By Proposition 4, it follows that

$\{a^*(B)\}^{-1}a(B) = cB^\ell$ for some constant c and some integer $\ell \in \mathbb{Z}$. Thus $a(B) = ca^*(B)B^\ell$ and the proof is complete. \square

Proof of Proposition 6

Similarly to (5), we have for $u, v \in \mathbb{R}$,

$$\log |\Psi^*(u)| + \log |\Psi^*(v)| = \sum_{k=1}^K \sum_{j=-\infty}^{\infty} \log |\Psi(ua_{k,j} + va_{k,j+h})|, \quad (\text{A8})$$

by the independence between $\epsilon_{k,t}$ and $\epsilon_{\ell,t'}$ for any $(k,t) \neq (\ell,t')$. By the arguments leading to (A4) we thus have

$$-|u|^{\alpha_*} L_1^* \left(\frac{1}{|u|} \right) = - \sum_{k=1}^K \sum_{j=-\infty}^{\infty} \sigma_k^{\alpha_k} |a_{k,j}|^{\alpha_k} |u|^{\alpha_k}, \quad (\text{A9})$$

in the neighborhood of 0. Thus $\alpha_k = \alpha_* =: \alpha$ for all k , and L_1^* is constant in the neighborhood of 0,

$$L_1^* \left(\frac{1}{|u|} \right) = \sum_{k=1}^K \sum_{j=-\infty}^{\infty} \sigma_k^\alpha |a_{k,j}|^\alpha.$$

Now by taking $v = tu$ in the neighborhood of $u = 0$, we find that, similarly to (A2),

$$\sum_{k=1}^K \sum_{j=-\infty}^{\infty} \sigma_k^\alpha |a_{k,j} + ta_{k,j+h}|^\alpha = (1 + |t|^\alpha) \sum_{k=1}^K \sum_{j=-\infty}^{\infty} \sigma_k^\alpha |a_{k,j}|^\alpha, \quad \forall t \in \mathbb{R}. \quad (\text{A10})$$

Now we can use the following elementary lemma.

LEMMA 2. *For $\alpha > 0$, let $f_\gamma : t \in \mathbb{R} \mapsto f_\gamma(t) = |1 + \gamma t|^\alpha$ and let $f_* : t \in \mathbb{R} \mapsto f_*(t) = |t|^\alpha$. Let $\Gamma = \{\gamma_i\}_{i \in \mathcal{I}}$ a family of distinct real numbers, with $\mathcal{I} \subset \mathbb{Z}$. The family of functions $\{f_*, f_\gamma; \gamma \in \Gamma\}$ is linearly independent.*

It follows from (A10) and Lemma 2 that for any (k, j) either $a_{k,j}$ or $a_{k,j+h}$ is equal to zero. Thus, for any k , the set $\{a_{k,j}, j \in \mathbb{Z}\}$ reduces to a singleton. The conclusion follows. \square

Proof of Proposition 7

By the independence between $\epsilon_{k,t}^*$ and $\epsilon_{\ell,t'}^*$ in the one hand, and $\epsilon_{k,t}$ and $\epsilon_{\ell,t'}$ in the other hand, for any $(k, t) \neq (\ell, t')$, we get from the i -th equation of (13)

$$\sum_{k=1}^K \sum_{j=-\infty}^{\infty} \log |\Psi^*(ua_{i,k,j}^* + va_{i,k,j+h}^*)| = \sum_{k=1}^K \sum_{j=-\infty}^{\infty} \log |\Psi(ua_{i,k,j} + va_{i,k,j+h})|.$$

Thus, by arguments already used,

$$\sum_{k=1}^K \sum_{j=-\infty}^{\infty} (\sigma_k^*)^{\alpha_k} |a_{i,k,j}^*|^{\alpha_k} |u|^{\alpha_k} = \sum_{k=1}^K \sum_{j=-\infty}^{\infty} \sigma_k^{\alpha_k} |a_{i,k,j}|^{\alpha_k} |u|^{\alpha_k},$$

in the neighborhood of 0. It follows that $\{\alpha_k; k = 1, \dots, K\} = \{\alpha_k^*; k = 1, \dots, K\}$ and thus $\alpha_1 = \alpha_1^* < \dots < \alpha_K = \alpha_K^*$ and $\sum_{j=-\infty}^{\infty} (\sigma_k^*)^{\alpha_k} |a_{i,k,j}^*|^{\alpha_k} |u|^{\alpha_k} = \sum_{j=-\infty}^{\infty} \sigma_k^{\alpha_k} |a_{i,k,j}|^{\alpha_k} |u|^{\alpha_k}$ for $k = 1, \dots, K$. It is thus not restrictive to assume $\sigma_k = \sigma_k^*$ and $\sum_{j=-\infty}^{\infty} |a_{i,k,j}^*|^{\alpha_k} |u|^{\alpha_k} = \sum_{j=-\infty}^{\infty} |a_{i,k,j}|^{\alpha_k} |u|^{\alpha_k}$. Now,

similarly to (A10),

$$\sum_{k=1}^K \sum_{j=-\infty}^{\infty} (\sigma_k^*)^{\alpha_k} |a_{i,k,j}^* + ta_{i,k,j+h}^*|^{\alpha_k} = \sum_{k=1}^K \sum_{j=-\infty}^{\infty} \sigma_k^{\alpha_k} |a_{i,k,j} + ta_{i,k,j+h}|^{\alpha_k}, \quad \forall t \in \mathbb{R},$$

which, using the fact that the powers α_k are different, entails, for $k = 1, \dots, K$,

$$\sum_{j=-\infty}^{\infty} |a_{i,k,j}^* + ta_{i,k,j+h}^*|^{\alpha_k} = \sum_{j=-\infty}^{\infty} |a_{i,k,j} + ta_{i,k,j+h}|^{\alpha_k}, \quad \forall t \in \mathbb{R}. \quad (\text{A11})$$

Let (η_t) be an i.i.d. sequence of symmetric α_k -stable distributed random variables, $\eta_t \sim \mathcal{S}(\alpha_k, 0, 1, 0)$. The left hand side of (A11) characterizes the distribution of the process $(a_{i,k}^*(B)\eta_t)$. It follows that $a_{i,k}^*(B)\eta_t \stackrel{d}{=} \frac{1}{\sigma_k} a_{i,k}(B)\epsilon_{k,t}$. Therefore, we have $\eta_t \stackrel{d}{=} \frac{1}{\sigma_k} \{a_{i,k}^*(B)\}^{-1} a_{i,k}(B)\epsilon_{k,t}$ and we can conclude as in the proof of Proposition 5. \square

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