Estimating the Codifference Function
of Linear Time Series Models with Infinite Variance

Dedi Rosadi* and Manfred Deistler

Institute for Mathematical Methods in Economics
Research Unit Econometrics and System Theory (EOS)
Vienna University of Technology, Argentinierstrasse 8, A1040 Vienna, Austria
E-mail: `manfred.deistler,dedi.rosadi@tuwien.ac.at`

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Abstract

We consider the *codifference* and the *normalized codifference* function as dependence measures for stationary processes. Based on the empirical characteristic function, we propose estimators of the *codifference* and the *normalized codifference* function. We show the consistency of the proposed estimators, where the underlying model is an ARMA with symmetric α-stable innovations, $0 < \alpha \leq 2$. In addition, we derive their limiting distribution. Finally, we present a simulation study showing the dependence of the estimator on certain design parameter.

Keywords: ARMA, Infinite Variance, Codifference, Empirical Characteristic Function

1 Introduction

In many cases, assumption of normality for the observations seems to be reasonable. On the other hand, in the number of applications, such as, signal processing, telecommunications, finance, physics and chemistry, the leptokurtic distribution, i.e., the distribution which is heavy-tailed and peaked around the center, seems to be more appropriate (e.g., Rachev and Mittnik, 2000; Nikias and Shao, 1995). An important class of distributions in this context is the stable distributions, which is a flexible class for data modelling and contains normal distributions as its special case. The importance of this class of distributions is strongly supported by generalized central limit theorems, which indicate that the stable distribution is the only possible limiting distribution for the normed sum of independent and identically distributed random variables. For information on the stable distribution, the reader is referred to, e.g., Samorodnitsky and Taqqu (1994).

In this paper, we consider univariate strictly stationary linear processes $\{X_t | t \in \mathbb{Z}\}$, where $\mathbb{Z}$ denotes the integers, given by

$$X_t = \sum_{j=0}^{\infty} c_j \epsilon_{t-j} \tag{1}$$

where the following holds

(C1). The coefficients c_j's are real-valued and satisfying $|c_j| < cQ^{-j}$ for some $c > 0$, $Q > 1$

(C2). ϵ_t is i.i.d. symmetric α stable ($S\alpha S$) distributed, i.e., ϵ_t has a characteristic function of the form

$$E \exp(is\epsilon_t) = \exp(-\sigma^\alpha |s|^\alpha) \tag{2}$$

where α denotes *the index of stability* $(0 < \alpha \leq 2)$, and $\sigma \geq 0$ denotes *the scale parameter*.

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*Permanent address: Department of Mathematics, Gadjah Mada University, Sekip Utara, Yogyakarta, Indonesia.

Under conditions C1 and C2, $\sum_{j=0}^{\infty} |c_j|^\alpha < \infty$, and the infinite sum (1) is well defined in the sense of a.s. convergence. Moreover, under the assumption C2, the process $\{X_t\}$ will be a strictly stationary $S\alpha S$ process with the same index of stability α but the scale parameter $\sigma_X = \sigma(\sum_{j=0}^{\infty} |c_j|^\alpha)^{1/\alpha}$ (Samorodnitsky and Taqqu, 1994, Theorem 7.12.2). Furthermore, the $ARMA(p,q)$ process for arbitrary (p,q),

$$\Phi(z)X_t = \Theta(z)\epsilon_t \tag{3}$$

for $t \in \mathbb{Z}$, where the polynomials Φ and Θ are given as

$$\Phi(z) = 1 - \phi_1 z - \phi_2 z^2 - \cdots - \phi_p z^p,$$
$$\Theta(z) = \theta_0 + \theta_1 z + \theta_2 z^2 + \cdots + \theta_q z^q,$$

with the real coefficients $\phi_1, \phi_2, \ldots, \phi_p$ and $\theta_1, \ldots, \theta_q$, and $\theta_0 = 1$, has unique stationary solution of the form (1) which fulfils condition C1 and C2 if and only if the polynomial $\Phi(z)$ has no roots in the closed unit disk $\{z : |z| \leq 1\}$. The polynomials $\Theta(z)$ and $\Phi(z)$ are assumed have no common roots, where z denotes backward-shift operator (here $z(X_t) = (X_{t-1})$) as well as the complex variable (Samorodnitsky and Taqqu, 1994, Theorem 7.12.2).

Here notice that if $\alpha = 2$, then ϵ_t is i.i.d. Gaussian with $\text{var}(\epsilon_t) = 2\sigma^2$. When $\alpha < 2$, $E|\epsilon_t|^p = \infty$ for $p \geq \alpha$ and $E|\epsilon_t|^p < \infty$ for $0 < p < \alpha$. Thus, the second moments of ϵ_t (and X_t) exist only for $\alpha = 2$, and for $\alpha < 2$, one can not use the covariance function $\gamma_n = E(X_n X_0)$ to describe dependence structure of the process $\{X_t\}$. Some generalizations of the autocovariance function as dependence measures of stationary process with infinite variance have been proposed in the literature, i.e., the *autocovariation* (see, e.g., Samorodnitsky and Taqqu, 1994), the *codifference* function (e.g., Kokoszka and Taqqu, 1994; Samorodnitsky and Taqqu, 1994) and the *dynamical function* (Janicki and Weron, 1994). In this paper, we consider the codifference function and analyze the properties of its estimator, both by an analytical investigation and the simulation study.

The rest of this paper is organized as follows. In section two, we present the main results of this paper. In this section, we give the definition of the codifference and the normalized codifference function, and also propose their estimator. Furthermore, for a class of linear processes, we show consistency of the sample codifference function, and further establish the limiting distribution of the proposed estimator. In section three, we present several simulation studies for the estimation of the normalized codifference function of pure moving average processes. Last section concludes.

2 Dependence structure of linear time series model with infinite Variance

2.1 Definition of the codifference function and its estimator

We consider the codifference function as proposed in Kokoszka and Taqqu (1994) and Yang *et al.* (2001)

$$\tau(k) = \tau(s, -s; k) = -\ln \mathrm{E}\exp(is(X_{t+k} - X_t)) + \ln \mathrm{E}\exp(isX_{t+k}) + \ln \mathrm{E}\exp(-isX_t) \tag{4}$$

where $s \in \mathbb{R}$ and $k \in \mathbb{Z}$. Because the characteristic function always exists, the codifference function requires no moment conditions for the original process $\{X_t\}$. In the Gaussian case, the codifference function is proportional to the covariance function, i.e., $\tau(s, -s; k) = -s^2\gamma(k)$, where $\gamma(\cdot)$ denotes the covariance function of the stationary process $\{X_t\}$. Moreover, by defining the normalized codifference function $I(k)$ as

$$I(k) = \frac{\tau(k)}{\tau(0)} = \frac{-\tau(k)}{-\tau(0)} \tag{5}$$

one directly obtains $I(k) = \rho(k)$ in the Gaussian case, where $\rho(k)$ denotes the correlation function.

Note that in general $\tau(-k) = \tau(k)^*$, where $\tau(k)^*$ denotes the conjugate of $\tau(k)$. For symmetric stationary process, $\tau(-k) = \tau(k)$ holds. In particular, under the assumptions C1 and C2, we obtain

that the codifference function $\tau(k)$ of the linear process (1) is of the form (see Kokoszka and Taqqu (1994))

$$\tau(k) = \sigma^\alpha |s|^\alpha \left[\sum_{j=0}^{\infty} (|(c_{j+k} - c_j)|^\alpha - |c_{j+k}|^\alpha - |-c_j|^\alpha) \right], k \geq 0 \qquad (6)$$

Note that under conditions C1 and C2, the normalized codifference function (5) is independent with the choice of s (i.e., for given α, it is equal to $I(1, -1; k) = \tau(1, -1; k)/\tau(-1, 1; 0)$ for any choice of s). For $S\alpha S$ process, $\tau(1, -1; k)$ is coincided with the codifference function $u(k)$ as given in Samorodnitsky and Taqqu (1994), eq. (4.7.1). For given strictly stationary $S\alpha S$ process X_t, $u(k)$ is defined as

$$u(k) = 2(\sigma_{X_t})^\alpha - (\sigma_{X_{t+k} - X_t})^\alpha \qquad (7)$$

where σ_Z and σ_{Y-Z} denote the scale parameters of Z and $Y - Z$, respectively.

Notice that if the $S\alpha S$ stationary process X_t is independent, then for $k \neq 0$, $u(k) = 0$, and clearly, $\tau(k) = 0$ for all s. Conversely, if $u(k) = 0, k \neq 0$ and $0 < \alpha < 1$, then X_t independent. When $1 \leq \alpha < 2$, $u(k) = 0$ does not imply that X_{t+k} and X_t are independent (Samorodnitsky and Taqqu, 1994).

Using Property 2.10.5 in Samorodnitsky and Taqqu (1994), we obtain that the normalized codifference $I(k)$ has the following property

$$0 \leq I(k) \leq 1 \text{ if } 0 < \alpha \leq 1 \qquad (8)$$

$$1 - 2^{\alpha - 1} \leq I(k) \leq 1 \text{ if } 1 \leq \alpha \leq 2 \qquad (9)$$

When $\alpha = 2$, (9) is equal to $-1 \leq \rho(k) \leq 1$, where $\rho(\cdot)$ denotes the autocorrelation function. Further theoretical properties of the codifference function were studied in Kokoszka and Taqqu (1994), Samorodnitsky and Taqqu (1994), Nowicka (1997) and Nowicka and Weron (1997).

As the codifference function is defined via characteristic functions (cf), it can be estimated by empirical characteristic functions (ecf)(see, e.g., Yu, 2004, for a review on ecf). Given a sample $X_1, X_2, \ldots, X_N$, an estimator for the codifference function at lag $k \in \mathbb{Z}$ can be defined as ($s \in \mathbb{R}$)

$$\hat{\tau}(s, -s; k) = \sqrt{(N-k)/N} \times [-\ln \phi(s, -s; k) + \ln \phi(s, 0; k) + \ln \phi(0, -s; k)] \qquad (10)$$

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

$$\phi(u, v; k) = \begin{cases} (N-k)^{-1} \sum_{t=1}^{N-k} \exp(i(uX_{t+k} + vX_t)) & \text{when } k \geq 0 \\ (N+k)^{-1} \sum_{t=1}^{N+k} \exp(i(uX_{t-k} + vX_t)) & \text{when } k < 0 \end{cases} \qquad (11)$$

Accordingly, $\hat{I}(s, -s; k) = \frac{\hat{\tau}(s, -s; k)}{\hat{\tau}(s, -s; 0)}$ can be used as the estimator of the normalized codifference $I(k)$. Here we consider a discrete estimation procedure, i.e., we evaluate the codifference function at r points $s_1 < \cdots < s_r$, for $s_i \in \mathbb{R}, s_i \neq 0, i = 1, \ldots, r$. In what follows, we denote the vectors $\mathbf{s} = \{s_1, \ldots, s_r\}$,

$$\hat{\tau}(\mathbf{s}, k) = [\hat{\tau}(s_1, -s_1; k), \hat{\tau}(s_2, -s_2; k), \ldots, \hat{\tau}(s_r, -s_r; k)]^T$$

and

$$\hat{I}(\mathbf{s}, k) = [\hat{I}(s_1, -s_1; k), \hat{I}(s_2, -s_2; k), \ldots, \hat{I}(s_r, -s_r; k)]^T$$

Note that one can replace the factor $\sqrt{(N-k)/N}$ in (10) by unity, and also the divisor $(N-k)$ in (11) by N without changing the asymptotic properties of the estimator, however, the choices in (10) and (11) will give a better finite sample performance than the alternative. Note that $\hat{\tau}(-k) = \hat{\tau}(k)$ such one can restrict the analysis to the case of $k \geq 0$. Two similar estimators for the codifference function have recently been proposed in Yang et $al.$ (2001) and in Hong (1999).

2.2 The asymptotic properties of the estimator

The asymptotic properties of the codifference estimator are summarized in the following theorems.

Theorem 2.1 *Let $X_t, t \in \mathbb{Z}$ be the stationary linear process (1) satisfying conditions C1 and C2. For $\mathbf{s} \in \mathbb{R}, s_i \neq 0, i = 1, \ldots, r$, its codifference estimator $\hat{\tau}(\mathbf{s}, k)$ and the sample normalized codifference $\hat{I}(\mathbf{s}, k)$ are (weakly) consistent estimators for $\tau(\mathbf{s}, k), k \in \{0, 1, 2, \ldots\}$ and $I(\mathbf{s}, k) = I(k)$, respectively.*

The proof is given in Appendix A.

The asymptotic distribution of the sample codifference function (and the sample normalized codifference function) of the linear process (1) can be derived using the central limit theorem for empirical characteristic function (Hesse, 1990, Theorem 1 and Remark 2.6.). For convenience, we split $\hat{\tau}$ into its real and imaginary parts. We write

$$\mathrm{Re}\,\hat{\tau}(\mathbf{s}, k) = [\mathrm{Re}\,\hat{\tau}(s_1, -s_1; k), \mathrm{Re}\,\hat{\tau}(s_2, -s_2; k), \ldots, \mathrm{Re}\,\hat{\tau}(s_r, -s_r; k)]^T$$

and

$$\mathrm{Im}\,\hat{\tau}(\mathbf{s}, k) = [\mathrm{Im}\,\hat{\tau}(s_1, -s_1; k), \mathrm{Im}\,\hat{\tau}(s_2, -s_2; k), \ldots, \mathrm{Im}\,\hat{\tau}(s_r, -s_r; k)]^T$$

Here, $\mathrm{Re}(z)$ and $\mathrm{Im}(z), z \in \mathbb{C}$ denote the real and imaginary parts of z. As $\hat{\tau}(s, -s, 0)$ by definition is a real function, we therefore obtain

$$\mathrm{Re}\,\hat{I}(\mathbf{s}, k) = \begin{bmatrix} \mathrm{Re}\,\hat{\tau}(s_1, -s_1; k)/\hat{\tau}(s_1, -s_1; 0) \\ \mathrm{Re}\,\hat{\tau}(s_2, -s_2; k)/\hat{\tau}(s_2, -s_2; 0) \\ \vdots \\ \mathrm{Re}\,\hat{\tau}(s_r, -s_r; k)/\hat{\tau}(s_r, -s_r; 0) \end{bmatrix}, \quad \mathrm{Im}\,\hat{I}(\mathbf{s}, k) = \begin{bmatrix} \mathrm{Im}\,\hat{\tau}(s_1, -s_1; k)/\hat{\tau}(s_1, -s_1; 0) \\ \mathrm{Im}\,\hat{\tau}(s_2, -s_2; k)/\hat{\tau}(s_2, -s_2; 0) \\ \vdots \\ \mathrm{Im}\,\hat{\tau}(s_r, -s_r; k)/\hat{\tau}(s_r, -s_r; 0) \end{bmatrix} \tag{12}$$

In the following theorem, a result regarding the asymptotic distribution of the sample normalized codifference is given. The proof is given in Appendix B.

Theorem 2.2 *Let $X_t, t \in \mathbb{Z}$ be the stationary linear process (1), satisfying conditions C1 and C2. Then for $h \in \{1, 2, \ldots\}$,*

$$\left[\begin{pmatrix} \mathrm{Re}\,\hat{I}(\mathbf{s}, 1) \\ \mathrm{Im}\,\hat{I}(\mathbf{s}, 1) \end{pmatrix}, \begin{pmatrix} \mathrm{Re}\,\hat{I}(\mathbf{s}, 2) \\ \mathrm{Im}\,\hat{I}(\mathbf{s}, 2) \end{pmatrix}, \ldots, \begin{pmatrix} \mathrm{Re}\,\hat{I}(\mathbf{s}, h) \\ \mathrm{Im}\,\hat{I}(\mathbf{s}, h) \end{pmatrix} \right]^T$$

is

$$AN\left(\left[\begin{pmatrix} I(1) \\ 0 \end{pmatrix}, \begin{pmatrix} I(2) \\ 0 \end{pmatrix}, \ldots, \begin{pmatrix} I(h) \\ 0 \end{pmatrix} \right]^T, N^{-1}\mathbf{W} \right) \tag{13}$$

The matrix variance-covariance $\mathbf{W}$ is given in (43).

Applying this theorem, we obtain the following corollary. The proof is given in Appendix C.

Corollary 2.3 *Let $X_t, t \in \mathbb{Z}$ be an i.i.d. sequence satisfying the condition C2. Then for $k \in \{1, 2, \ldots\}$,*

$$\mathrm{Re}\,\hat{I}(\mathbf{s}, k) \text{ is } AN(0, N^{-1}\mathbf{W}_1) \tag{14}$$

and

$$\mathrm{Im}\,\hat{I}(\mathbf{s}, k) \text{ is } AN(0, N^{-1}\mathbf{W}_2) \tag{15}$$

where the (i, j)-th elements of matrix $\mathbf{W}_1$ and $\mathbf{W}_2$ are,

$$W_1(i, j) = \frac{f_{ij}}{g_{ij}} \text{ and } W_2(i, j) = \frac{h_{ij}}{g_{ij}}, i, j = 1, \ldots, r \tag{16}$$

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with

$$f_{ij} = e^{\sigma^\alpha(|s_i|^\alpha + |s_j|^\alpha - |s_i - s_j|^\alpha)} \left\{ \tfrac{1}{2} e^{\sigma^\alpha(|s_i|^\alpha + |s_j|^\alpha - |s_i - s_j|^\alpha)} - 1 \right\}$$
$$+ e^{\sigma^\alpha(|s_i|^\alpha + |s_j|^\alpha - |s_i + s_j|^\alpha)} \left\{ \tfrac{1}{2} e^{\sigma^\alpha(|s_i|^\alpha + |s_j|^\alpha - |s_i + s_j|^\alpha)} - 1 \right\} + 1$$

$$h_{ij} = e^{\sigma^\alpha(|s_i|^\alpha + |s_j|^\alpha - |s_i - s_j|^\alpha)} \left\{ \tfrac{1}{2} e^{\sigma^\alpha(|s_i|^\alpha + |s_j|^\alpha - |s_i - s_j|^\alpha)} - 1 \right\}$$
$$+ e^{\sigma^\alpha(|s_i|^\alpha + |s_j|^\alpha - |s_i + s_j|^\alpha)} \left\{ 1 - \tfrac{1}{2} e^{\sigma^\alpha(|s_i|^\alpha + |s_j|^\alpha - |s_i + s_j|^\alpha)} \right\}$$

and

$$g_{ij} = 4\sigma^{2\alpha} |s_i|^\alpha |s_j|^\alpha$$

3 Simulation evidence

In this section, we present a simulation study for investigating the small sample properties of the sample normalized codifference function.

3.1 Practical considerations

Before we proceed, we make a remark about the sample and the population codifference function. From (6) and the fact that all c_j's are real, we have that the codifference function of the models we consider here is a real-valued function but the estimator (10) is complex. Therefore, one possibility is to use only the real part of the estimator. Because in practice we are working with a finite sample, the imaginary part of the estimator is still present, but will vanish asymptotically.

In what follows, we are only working with the estimator of normalized codifference $\hat{I}(k)$. Hence, we note that unlike the true normalized codifference (5), from (10) and Corollary 2.3, one can see that the sample normalized codifference function and its limiting variance depend on $\mathbf{s} = \{s_1, \ldots, s_r\}$. Apparently $\hat{I}(\cdot)$ is defined for all $\mathbf{s} > \mathbf{0}$, and from Theorem 2.1, we know that it is a consistent estimator. However, in a finite sample, the convergence of the estimator to the population values depends on the choice of $\mathbf{s}$. Therefore, for estimation, $\mathbf{s}$ is a design parameter which has to be chosen appropriately. In other words, $\hat{I}(\cdot)$ should be calculated from those values of $\mathbf{s}$ which gives the most accurate estimates of the true function $I(\cdot)$.

To be more precise, in practice the number of grid points r and more importantly, the location of $s_1, \ldots, s_r$, have to be chosen. As the normalized codifference function is defined based on the *ecf*, we can apply here the known results for *ecf*. For a fixed r, Koutrouvelis (1980) and Kogon and Williams (1998) showed that for calculating *ecf*, the location of the grid points $s_1, \ldots, s_r$, should be chosen close, but not equal to zero. It has been shown in this case, the *ecf* will most accurately estimate the characteristic function. This particular choice of grid points seems to be reasonable for calculating $\hat{I}(\cdot)$, as for instance, can be shown in Figure 1. To determine the location of s_i's, we suggest to plot $\operatorname{Re} \hat{I}(k)$ within the interval $0.01 \le s \le 2$, for some values of lag $k > 0$. These graphs will show the interval of s around zero which has relatively small bias, as the best location for evaluating the estimator. The best choice for the interval of s clearly depends on the data itself and in general also on the lag k. However, we suggest to use $s_a = 0.01$ as the left bound of the interval, where the best choice for the right bound can be determined from the graphs, i.e., as the threshold of $s = s_b$ where the graphs $\operatorname{Re} \hat{I}(k)$, for some lag k, are still relatively flat. The individual choices for s_i's can be chosen in one of two ways:

1. If we wish to use equal spacing of s_i's, we can set $\mathbf{s} = \{s_1 = 0.01, 0.01 + i\frac{s_b - s_a}{r-1}, s_b\}, i = 1, \ldots, r - 2$. To obtain r, we can minimize the determinant of covariance matrix in (13). Unfortunately, the covariance matrix in (13) depends on the unknown parameters c_j's, α, σ and the distance between s_i's. One possibility is to replace (13) with its consistent estimate, however here we consider a different approach for choosing the s_i's, i.e., with the help of

the trajectories of the estimator, for instance, as given in Figure 1. Figure 1 indicates that these trajectories will depend on the sample size N and, more strongly, on α. From our numerical studies, we observe that for given α, the behavior of these trajectories is typical for arbitrary lag k. For $\alpha = 2$, $\hat{I}(k)$ is relatively smooth, where for $\alpha < 2$, $\hat{I}(k)$ has an erratic behavior, and this behavior will be stronger when α is goes farther away from 2. This result suggests that when $\alpha = 2$, there is no benefit by choosing the s_i's very close to each other, and conversely for $\alpha < 2$. This conjecture can be checked in i.i.d. case using the determinant of the covariance matrix (14). Here, we propose to use s_i's with distance $d = 0.01$ for $\alpha \leq 1$, $0.01 < d \leq 0.05$ for $1 < \alpha \leq 1.5$, $0.05 < d \leq 0.1$ for $1.5 < \alpha < 2$ and $d = 0.1$ or larger for $\alpha = 2$. Especially in i.i.d. case we can show that these choices are sufficient, in the sense that for given α, choosing a smaller distance between grid points will not significantly decrease the determinant of the covariance matrix (14). Notice that in practice it is not necessary to know α. As the erratic behavior of the estimator is typical for given α, we can observe this property using the plot of $\operatorname{Re} \hat{I}(k)$ within an interval near zero, for some values of lag $k > 0$.

2. If equal spacing is not considered, when r has been fixed, the choice of s_i's can be chosen using the determinant of covariance matrix in (13). However, here we can use a similar consideration as above, i.e., we choose the s_i's sufficiently close, depending on the erratic behavior of the estimates $\operatorname{Re} \hat{I}(\cdot)$.

The last thing to consider is the number of points r. For $r \geq 1$, the final estimate $\hat{I}(k)$ can be defined as the weighted average of the estimates at the grid points $s_1, \ldots, s_r$, i.e., we define $\hat{I}(k) = \sum_{i=1}^{r} w_i \hat{I}(s_i, -s_i; k)$ with $\sum_{i=1}^{r} w_i = 1$. For instance, we can use a simple average with $w_i = 1/r$ or a negative exponentially weighted average with $w_i = \exp(-s_i^2) / \sum_{j=1}^{r} \exp(-s_j^2)$. In i.i.d. case, we obtain that by averaging the estimates at different points, the asymptotic variance of the estimator will be smaller or equal to the variance of the estimator obtained at single point, which can be seen directly from Figure 2. Figure 2 shows that for $\alpha = 2$, there is no difference in terms of the asymptotic variance between estimating $\hat{I}(\cdot)$ either at a single point or at more points, whereas for $\alpha < 2$, the difference is significant, especially when α is small. Furthermore, Figure 2 also shows that the smaller α is, the smaller the covariance between the grid points. For the finite sample case, this fact agrees with the typical erratic behavior of the plot $\operatorname{Re} \hat{I}(\cdot)$ of some non i.i.d. samples, for instance, as shown in Figure 1. Based on these results and from our numerical studies, in the finite sample case, we suggest the choice of the number of grid points r as follows. For $\alpha = 2$ (i.e., for smooth graphs of $\operatorname{Re} \hat{I}(k)$, $k > 0$), we observe that $r = 1$ is sufficient, whereas for $\alpha < 2$ (i.e., for erratic graphs of $\operatorname{Re} \hat{I}(k)$, $k > 0$), at least two points should be chosen, and more points are required when α is farther away from 2 (i.e., for more erratic graphs of $\operatorname{Re} \hat{I}(k)$, $k > 0$). It is important to note that from our simulation experience, the accuracy of the estimator is more sensitive to the location of grid points than to the number of points r.

In the following subsection, we will investigate the choice of grids points through monte carlo simulations.

3.2 Simulation results

To investigate the proposed choice of s and also the finite sample behavior of the estimator, we run several monte-carlo simulations using $R/GNU\text{-}S$ version 1.7.0 (R Development Core Team, 2004) (available on the Web at `http://www.r-project.org`), where we use function $rstable$ in the extension package `stable` (available on the Web at `http://alpha.luc.ac.be/~jlindsey/rcode.html`), to generate the unit symmetric α stable innovations (based on method presented in Chambers *et al.*, 1976) and function $arima.sim$ in the package `stats` to generate $X_t = \epsilon_t + c_1 \epsilon_{t-1} + c_2 \epsilon_{t-2}$ processes where (c_1, c_2) are

$$I.\,(2, 1.111) \quad II.\,(-1, 0.5) \quad III.\,(0.55, 0.05) \quad IV.\,(-0.4, 0.7) \tag{17}$$

and from now on we refer to these as experiment I - experiment IV, respectively.

In Gaussian framework, models in experiments I, II and III were examined in Bhansali (1983). The roots of the polynomial $1 + c_1 z + c_2 z^2 = 0$ are as follows. In experiment I, the roots are

$-0.9 \pm 0.3i$, close to the invertibility region. In experiment II and IV, the roots are $0.5 \pm 0.5i$, and $-0.2857 \pm 0.247i$, so the absolute values of the roots are 0.71, and 0.378, respectively. In these experiments, the models have similar roots properties which are neither too close to 1 nor to 0. In experiment III, the roots are real-valued, equal to -0.435 and -0.115, one close to 0.5 and the other close to 0.

For $\alpha = 2$, the true values of the normalized codifference (equal to the correlation function) at lag k, $(I(1), I(2)) = (\rho(1), \rho(2))$ in experiment I-IV are:

$$(I). (0.677, 0.178), \quad (II). (-0.667, 0.222), \quad (III). (0.443, 0.038), \quad (IV). (-0.412, 0.424)$$

In experiment I and II, the values of I(k) are closer to 1 at lag 1 and not too close to 0 at lag 2 while for experiment III, at lag 1 close to 0.5 but almost 0 at lag 2. For the last experiment, at lag 1 it is negative but it is positive at lag 2, with absolute values near 0.5.

In order to see the performance of the estimator, we simulate the time series in experiment I - IV for several values of α with $\sigma = 1$ and two sample sizes, the the "small" one is $N = 100$ and the "large" one is $N = 1000$. All experiments are replicated $T = 1000$ times. The estimates of the normalized codifference are calculated for lags 1 till lag 10. Figure 1 suggests that in the interval $0.01 \leq s \leq 0.5$, $\mathrm{Re}\,\hat{I}(\cdot)$ is relatively less biased, although the best interval of $\mathbf{s}$ depends on the index α. For checking the best location of $\mathbf{s}$ and also the choice of grid points, we choose several different sets of $\mathbf{s}_i = \{s_1, \ldots, s_r\}, i = 1, 2, \ldots, 28$. Here we consider the equidistant and non equidistant grid points. The complete listing of the choices is as follows : $\mathbf{s}_1 = \{0.01\}$, $\mathbf{s}_2 = \{0.1\}$, $\mathbf{s}_3 = \{0.2\}$, $\mathbf{s}_4 = \{0.3\}$, $\mathbf{s}_5 = \{0.5\}$, $\mathbf{s}_6 = \{1\}$, $\mathbf{s}_7 = \{0.01, 0.1\}$, $\mathbf{s}_8 = \{0.01, 0.2\}$, $\mathbf{s}_9 = \{0.01, 0.5\}$, $\mathbf{s}_{10} = \{0.01, 1\}$, $\mathbf{s}_{11} = \{0.1, 0.2\}$, $\mathbf{s}_{12} = \{0.1, 0.5\}$, $\mathbf{s}_{13} = \{0.1, 1\}$, $\mathbf{s}_{14} = \{0.5, 1\}$, $\mathbf{s}_{15} = \{0.01, 0.1, 0.2\}$, $\mathbf{s}_{16} = \{0.01, 0.1, 0.5\}$, $\mathbf{s}_{17} = \{0.01, 0.1, 1\}$, $\mathbf{s}_{18} = \{0.1, 0.2, 0.3\}$, $\mathbf{s}_{19} = \{0.1, 0.3, 0.5\}$, $\mathbf{s}_{20} = \{0.01, 0.5, 1\}$, $\mathbf{s}_{21} = \{0.1, 0.5, 1\}$, $\mathbf{s}_{22} = \{0.1, 0.2, 0.3, 0.4, 0.5\}$, $\mathbf{s}_{23} = \{0.1, 0.2, \ldots, 1\}$, $\mathbf{s}_{24} = \{0.01, 0.06, 0.11, 0.16, 0.21\}$, $\mathbf{s}_{25} = \{0.01, 0.02, \ldots, 0.2\}$, $\mathbf{s}_{26} = \{0.01, 0.02, \ldots, 0.1\}$, $\mathbf{s}_{27} = \{0.11, 0.12, \ldots, 0.2\}$ and $\mathbf{s}_{28} = \{0.5, 0.55, \ldots, 1\}$. For each choice of $\mathbf{s}_i$ in run h, the final estimates are calculated as the weighted value of estimates among the choices of grid points $s_{ij}, j = 1, \ldots, r_i$, denoted by $\mathrm{Re}\,\hat{I}(\cdot)_{ih} = \sum_{j=1}^{r_i} w_{ij}\,\mathrm{Re}\,\hat{I}(s_{ij}, -s_{ij}, \cdot)_h$, where $\mathrm{Re}\,\hat{I}(s_{ij}, -s_{ij}, \cdot)_h$ denotes the real part of the estimates of the normalized codifference in run h at certain lags, calculated at $s_{ij}, j = 1, \ldots, r_i$. Here we consider two methods for weighting the estimates, first we use a simple average of the estimates and the second, we use a negative exponential weighted average. To save the space, we only present the result of experiment 1, which is summarized in table 1, but the results in the other experiments are equivalent. In the table, we also record the best choices of $\mathbf{s}$, which are defined as the values of grid points which minimize the sum of mean absolute deviation (MAD) of estimates at lag 1 and lag 2, among all considered choices of grid points above. Here, MAD at lag k and for grid $\mathbf{s}_i$ is defined as $MAD_{ik} = \frac{1}{T} \sum_{h=1}^{T} | \mathrm{Re}\,\hat{I}(k)_{ih} - I(k)|$, $k = 1, 2$. For the sake of comparison, when $\alpha = 2$ we also record the estimates of sample (central) ACF.

From the result of simulation, as expected, we observe that the estimation accuracy will be improved when the sample size is increased. Furthermore, throughout the simulation studies, the results indicate that the accuracy of the estimates of the normalized codifference function depends on the choice of grid points $\mathbf{s}$, where the optimal choices of grids depend on the index α and the sample size N. When $\alpha = 2$, surprisingly that under suitable choices of grid points, we find in some cases, the sample normalized codifference can provide a better estimate (in terms of total MAD for the first two lags) to the true values (of the normalized codifference, which is equal to ACF) than the estimates given by the sample ACF. When $\alpha < 2$, it seems that there is a great benefit by evaluating $\mathrm{Re}\,\hat{I}(\cdot)$ at several points of s, that is $r > 1$, where under the appropriate choice of grid points $\mathbf{s}$, the performance of the weighting methods (the simple average and the exponential weight) are approximately the same. For MA(2) models, in all cases we consider here, we also find that the estimation accuracies are significantly better if we choose the grids point $s < 0.5$ than $s \geq 0.5$. In general, there is a benefit in terms of the estimation accuracy to include a point close to zero. We further observe that when $\alpha < 1.5$, the choice of equidistant grids with a distance between 0.01 and 0.05 seems to be adequate. For $\alpha \geq 1.5$, a distance 0.1 seems to be adequate, because a smaller grid distance does not really improve the accuracy of the estimates.

As a general conclusion, from this simulation studies, we may conclude that the optimal choice of grid points **s** will follow the lines of our proposed choice of grid points **s** as in Section 3.1.

4   Conclusion

In this paper, we propose estimators of the codifference and the normalized codifference function, where for the linear processes with geometrically bounded coefficients and $S\alpha S$ innovations, we established the asymptotic properties of the proposed estimators. Notice that unlike the ACF estimator, we obtain that there is no discontinuity in either the normalization or the limiting distribution of the proposed estimators when $\alpha \rightarrow 2$. Moreover, we note that unlike the sample ACF which has an unfamiliar limiting distribution when $\alpha < 2$ and relatively difficult to obtain the quantiles of the limit distribution, estimators of the codifference and the normalized codifference will be asymptotically normally distributed at the same rate as the sample ACF in the classical case although the asymptotic variance is different. We also present a simulation study to observe the small sample properties of the normalized codifference estimator.

In the practical situation, to obtain the estimates of the normalized codifference with a good accuracy for the real data, first we suggest to plot $\mathrm{Re}\,\hat{I}(k)$ within the interval $s \in [0.01, 2]$, for some values of lag $k > 0$. These graphs will show two important things. First, they suggest the interval of s near zero which has a small bias (i.e., the interval $0.01 < s < s_b$, s_b denotes the threshold point of s where the graphs $\mathrm{Re}\,\hat{I}(k)$, for some k, are still relatively flat), as the best location for evaluating $\hat{I}(\cdot)$. Secondly, it reveals the erratic behavior of the estimates. When the graphs are smooth, the choice of one point $s = 0.01$ is sufficient for estimating $\mathrm{Re}\,\hat{I}(k)$. If the graphs are erratic, at least two points are required and more points are better for more erratic graphs. If the equidistant points $s_1, \ldots, s_r$ are considered, when the graphs are highly erratic, we can use a small distance between points, e.g., 0.01, where for less erratic graphs, we can use a bigger distance, such as 0.05 or 0.1. If the non equidistant points are used, we should include one point close to zero in the choice of s_i's and the chosen points are sufficiently close to each other. Finally, we define the final estimate as the weighted average of the estimates at $s_1, \ldots, s_r$.

Notice that in this paper we have considered a method for calculating the codifference and the normalized codifference "directly" from the data. As an obvious alternative, once one knows the estimated parameters and the orders of the estimated models, one may directly estimate the codifference and the normalized codifference using equation (6). The methods for estimating the parameters of stable ARMA models have been reviewed in, e.g., Embrechts *et al.* (1997), Chapter 7. Notice that for small order MA and AR processes, the tail index α can be well estimated using a quantile based estimator (i.e McCulloch's method), see, e.g., Adler *et al.* (1998) for simulation evidences. In our opinion, for inference purpose, the "direct" estimation method is more preferable than estimating the codifference function via estimated parameters.

APPENDIX

A   Proof of Theorem 2.1

Before we give the lemmas which are necessary for the consistency proof of the codifference estimator, a related result from Kokoszka and Taqqu (1994) will be presented, which shows the codifference function in ARMA case is bounded by an exponentially decaying function just like the covariance function in the classical case. Kokoszka and Taqqu (1994) consider more general

definition of the codifference function (for $\theta_1, \theta_2 \in \mathbb{R}$)

$$\tau_G(\theta_1, \theta_2; k) = -\ln \mathrm{E} \exp(i(\theta_1 X_{t+k} + \theta_2 X_t)) + \ln \mathrm{E} \exp(i\theta_1 X_{t+k}) + \ln \mathrm{E} \exp(i\theta_2 X_t) \qquad (18)$$

but contain (4) as the special case ($\theta_1 = s, \theta_2 = -s$).

Theorem A.1 (Kokoszka and Taqqu (1994), Theorem 2.1.) *If the coefficients c_j's of the linear process (1), satisfying conditions C1 and $\{\epsilon_t\}$ satisfying C2 then ($\theta_1, \theta_2 \in \mathbb{R}$)*

$$\limsup_{k \to \infty} Q^{\alpha k} |\tau_G(k)| \le 2(1 - Q^\alpha)^{-1/\alpha} |\theta_1|^\alpha \ \text{for } 0 < \alpha \le 1 \qquad (19)$$

and

$$\limsup_{k \to \infty} Q^k |\tau_G(k)| \le \alpha \left(\sum_{j=0}^\infty |c_j|^\alpha \right)^{\frac{\alpha - 1}{\alpha}} (1 - Q^\alpha)^{-1/\alpha} |\theta_1| |\theta_2|^{\alpha - 1} \ \text{for } 1 < \alpha \le 2 \qquad (20)$$

To show consistency of the codifference estimator, the following two lemmas are necessary.

Lemma A.2 *Let $X_t, t \in \mathbb{Z}$ be the stationary linear process (1), satisfying conditions C1 and C2, and let $\Phi(s) = \mathrm{E} \exp(isX_t)$ denote its first order characteristic function. For $k \in \{0, 1, 2 \ldots\}$ and $s \in \mathbb{R}, s \ne 0$*

$$\ln \phi(s, k) = \ln \left((N - k)^{-1} \sum_{t=1}^{N-k} \exp(isX_t) \right)$$

is a consistent estimator of $\ln \Phi(s)$.

Proof. Let $y_t = \exp(isX_t)$. Apparently, the magnitude of y_t is equal to one, and therefore it is a second order stationary process. For the notation simplicity, instead of working with $\phi(s, k)$, we first show consistency of $\phi^*(s, k) = N^{-1} \sum_{t=1}^N \exp(isX_t)$. Here, $\phi^*(s, k)$ is an unbiased estimator for $\Phi(s) = E(y_t)$. To show the weak consistency of this estimator, we show that y_t is a mean ergodic process. A sufficient condition for y_t to be mean ergodic, i.e., $\phi^*(s, k) \to E(y_t)$ in the mean square sense, is that its covariance function tends to zero as time lags tends to ∞ (e.g., Theorem 7.1.1. in Brockwell and Davis, 1987). The covariance function of y_t at lag k can be expressed as

$$c(k) = |\Phi(s)|^2 \left(\frac{\mathrm{E}(\exp(is(X_{t+k} - X_t)))}{\mathrm{E}(\exp(isX_{t+k}))\mathrm{E}(\exp(-isX_t))} - 1 \right) = |\Phi(s)|^2 (\exp(-\tau(k)) - 1) \qquad (21)$$

From Theorem A.1, we see that $c(k) \to 0$ when $k \to \infty$ exponentially fast. As mean square convergence entails convergence in probability, $\phi^*(s, k) \xrightarrow{p} \Phi(s)$. Moreover, under assumptions C1 and C2, we have $\Phi(s) = \exp(- \sum_{j=0}^\infty \sigma^\alpha |sc_j|^\alpha)$, a real-valued function. Therefore we can conclude $\mathrm{Re}\, \phi^*(s, k) \xrightarrow{p} \mathrm{Re}\, \Phi(s) = \Phi(s)$ and $\mathrm{Im}\, \phi^*(s, k) \xrightarrow{p} \mathrm{Im}\, \Phi(s) = 0$.

By taking the principal value of $\ln(\cdot)$ function in the complex domain, we can see that $\ln(\cdot)$ is a continuous and well-defined function on $\mathbb{C}$ minus the negative real line. Because $|c_j| < cQ^{-j}$ for some $c > 0$, $Q > 1$, we conclude $\mathrm{Re}\, \Phi(s)$ always strictly greater than 0, which implies with the probability converging to 0, $\mathrm{Re}\, \phi^*(s, k)$ will be less than or equal to 0. Therefore, without loss of generality, we can restrict the definition of the real and imaginary parts of $\ln \phi^*(s, k)$ only on the right half plane where $\mathrm{Re}\, \phi^*(s, k) > 0$, and equal to 0 on the other case. From this consideration, we obtain $\mathrm{Re} \ln \phi^*(s, k) = \frac{1}{2} \ln((\mathrm{Re}\, \phi^*(s, k))^2 + (\mathrm{Im}\, \phi^*(s, k))^2)$ and $\mathrm{Im} \ln \phi^*(s, k) = \arctan(\frac{\mathrm{Im}\, \phi^*(s, k)}{\mathrm{Re}\, \phi^*(s, k)})$. From the continuity of the logarithm function in the considered domain, we can deduce that $\mathrm{Re} \ln \phi^*(s, k) \xrightarrow{p} \mathrm{Re} \ln \Phi(s) = \ln \Phi(s)$ and $\mathrm{Im} \ln \phi^*(s, k) = \arg \phi^*(s, k) \xrightarrow{p} 0$, when $N \to \infty$. In other words, we obtain $\ln \phi^*(s, k) \xrightarrow{p} \ln \Phi(s)$. To complete our proof, it is sufficient to show $\phi^*(s, k) - \phi(s, k) \xrightarrow{p} 0$. By assumption of the model, $\mathrm{Re}\, \Phi(s) > 0$, thus $E |\mathrm{Re}\, \phi^*(s, k) - \mathrm{Re}\, \phi(s, k)| < 2\frac{k}{N-k}$ and $E |\mathrm{Im}\, \phi^*(s, k) - \mathrm{Im}\, \phi(s, k)| < 2\frac{k}{N-k}$, and therefore we can conclude $\phi^*(s, k) - \phi(s, k) = o_p(1)$. ∎

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Lemma A.3 *Let $X_t, t \in \mathbb{Z}$ be the stationary linear process (1), satisfying conditions C1 and C2, and for $k \in \{0, 1, 2, \ldots\}$ and $s \in \mathbb{R}, s \neq 0$, let $\Phi(s, -s; k) = \mathrm{E}\exp(is(X_{t+k} - X_t))$ be its second-order characteristic function evaluated at $(s, -s)$. Then as $N \to \infty$*

$$\ln \phi(s, -s; k) \xrightarrow{p} \ln \Phi(s, -s; k)$$

where $\phi(s, -s; k)$ is as given in (11).

Proof. For the proof, we can proceed in a similar way as the previous lemma. For simplicity, instead of working with $\phi(s, -s; k)$, we first show the consistency of $\phi^*(s, -s; k) = N^{-1}\sum_{t=1}^{N} \exp is(X_{t+k} - X_t)$. A sufficient condition for y_t to be autocovariance ergodic (Proakis and Manolakis, 1996, p.A10), i.e., $\phi^*(s, -s, k) \to \Phi(s, -s; k)$, in the mean square sense is that

$$\mathrm{E}\exp(is(X_t - X_{t+k} - X_{t+n} + X_{t+n+k})) \to |\Phi(s, -s; k)|^2$$

as $n \to \infty$ where the index n denotes the lag of covariance among the sample autocovariance function. Hence, we have

$$\mathrm{E}\exp(is(X_t - X_{t+k} - X_{t+n} + X_{t+n+k}))$$
$$= |\Phi(s, -s; k)|^2 \frac{\mathrm{E}\exp(is((X_t - X_{t+k}) - (X_{t+n} - X_{t+n+k})))}{\mathrm{E}\exp(is(X_t - X_{t+k}))\mathrm{E}\exp(is(X_{t+n+k} - X_{t+n}))} = |\Phi(s, -s; k)|^2 \exp(-C_n)$$

where

$$C_n = -\ln \mathrm{E}\exp(is((X_t - X_{t+k}) - (X_{t+n} - X_{t+n+k}))) \tag{22}$$
$$+ \ln \mathrm{E}\exp(is(X_t - X_{t+k})) + \ln \mathrm{E}\exp(is(X_{t+n+k} - X_{t+n})) \tag{23}$$

Applying the similar technique as obtaining (6), one can write C_n as

$$C_n = \sigma^\alpha \left[\sum_{j=0}^{\infty} |s(c_j - c_{j+k} - c_{j+n} + c_{j+n+k})|^\alpha - |s(c_{j+n+k} - c_{j+n})|^\alpha - |s(c_j - c_{j+k})|^\alpha \right]$$

$$= \sigma^\alpha \left[\sum_{j=0}^{\infty} |s(k_j - k_{j+n})|^\alpha - |-sk_{j+n}|^\alpha - |sk_j|^\alpha \right]$$

where $k_j = c_j - c_{j+k}$. This expression is the codifference function $\tau_G(n)$ for coefficients k_j's and parameters $\theta_1 = -s$, $\theta_2 = s$. Because $|c_j| < cQ^{-j}$ for some $c > 0$, $Q > 1$, then $|k_j| < c_1 Q^{-j}$ for some $c_1 = 2c > 0$, $Q > 1$. Therefore, by (20) and (19), we can conclude that $\exp(-C_n)$ will converge to 1 exponentially fast. In other words, $\mathrm{E}\exp(is(X_t - X_{t+k} - X_{t+n} + X_{t+n+k})) \to |\Phi(s, -s; k)|^2$ for $n \to \infty$, and we obtain the mean square convergence of $\phi^*(s, -s; k)$ to $\Phi(s, -s; k)$ and therefore $\phi^*(s, -s; k) \xrightarrow{p} \Phi(s, -s; k)$. For the rest of the proof, we can proceed similarly to the proof of previous lemma, as we have $\Phi(s, -s; k) = \exp(-\sigma^\alpha(\sum_{j=0}^{k-1} |sc_j|^\alpha - \sum_{j=0}^{\infty} |s(c_{j+k} - c_j)|^\alpha))$ also a real-valued function, strictly greater than 0. ∎

Proof of Theorem 2.1. As for finite k and $N \to \infty$ we obtain $\sqrt{1 - k/N} \to 1$, then using the results in lemma A.2 and lemma A.3, we have as $N \to \infty$, for $i = 1, \ldots, r$

$$\hat{\tau}(s_i, -s_i; k) \xrightarrow{p} -\ln \Phi(s_i, -s_i; k) + \ln \Phi(s_i) + \ln \Phi(-s_i) = \tau(s_i, -s_i; k) \tag{24}$$

∎

B   The limit distribution of the sample codifference function

In this part, we will derive the asymptotic distribution of the sample codifference function of linear processes. The proof will be given as a series of propositions, where the main results are presented in Theorem B.4 and also the proof of Theorem 2.2 at the end of this part. The proof will follow closely an approach for obtaining the limiting distribution of the sample ACF in the classical case, e.g., Theorem 7.2.1 in Brockwell and Davis (1987).

For notational simplicity, instead of working with $\hat{\tau}(s_i, -s_i; k)$, $i = 1, \ldots, r$, in the following first we will consider the similar estimator $\hat{\tau}^*(s_i, -s_i; k)$,

$$\hat{\tau}^*(s_i, -s_i; k) = -\ln \phi^*(s_i, -s_i; k) + \ln \phi^*(s_i, 0; k) + \ln \phi^*(0, -s_i; k) \tag{25}$$

where $\phi^*(u, v; k) = N^{-1} \sum_{t=1}^{N} \exp(i(uX_{t+k} + vX_t))$, $u, v \in \mathbb{R}$. The required result will be presented in Theorem B.4.

Proposition B.1 *Let $X_t, t \in \mathbb{Z}$ be the stationary linear process (1), satisfying conditions C1 and C2. Then if $p \geq 0$ and $q \geq 0$,*

$$\lim_{N \to \infty} N \mathrm{cov}\left(\left(\begin{array}{c} \mathrm{Re}\, \hat{\tau}^*(\mathbf{s}, p) \\ \mathrm{Im}\, \hat{\tau}^*(\mathbf{s}, p) \end{array} \right), \left(\begin{array}{c} \mathrm{Re}\, \hat{\tau}^*(\mathbf{s}, q) \\ \mathrm{Im}\, \hat{\tau}^*(\mathbf{s}, q) \end{array} \right) \right) = \lambda \mathbf{L}_2^p \mathbf{V}_{pq} \mathbf{L}_2^q \lambda^T$$

where the matrices $\lambda, \mathbf{L}_2^k, k = p, q$ and $\mathbf{V}_{pq}$ are given in (27), (34) and (36) below. Here $\mathrm{cov}(X, Y)$ denotes the covariance between X and Y.

Proof. To obtain a complete variance-covariance structure of the estimator, we consider the following representation of $\hat{\tau}^*(\mathbf{s}, k)$

$$\left(\begin{array}{c} \mathrm{Re}\, \hat{\tau}^*(\mathbf{s}, k) \\ \mathrm{Im}\, \hat{\tau}^*(\mathbf{s}, k) \end{array} \right) = \left(\begin{array}{c} \mathrm{Re}\, \hat{\tau}^*(s_1, -s_1, k) \\ \mathrm{Re}\, \hat{\tau}^*(s_2, -s_2, k) \\ \vdots \\ \mathrm{Re}\, \hat{\tau}^*(s_r, -s_r, k) \\ \mathrm{Im}\, \hat{\tau}^*(s_1, -s_1, k) \\ \mathrm{Im}\, \hat{\tau}^*(s_2, -s_2, k) \\ \vdots \\ \mathrm{Im}\, \hat{\tau}^*(s_r, -s_r, k) \end{array} \right) = \lambda \left(\begin{array}{c} \mathbf{Y} \\ \mathbf{X} \end{array} \right) \tag{26}$$

where

$$\lambda = \left(\begin{array}{cc} \mathbf{I}_r \otimes \lambda_1 & \mathbf{0} \\ \mathbf{0} & \mathbf{I}_r \otimes \lambda_1 \end{array} \right) \tag{27}$$

$$\lambda_1 = \left(\begin{array}{ccc} 1 & 1 & -1 \end{array} \right)$$

and

$$\mathbf{Y} = \left(\begin{array}{c} \mathrm{Re}\ln Y_1^k \\ \mathrm{Re}\ln Y_2^k \\ \vdots \\ \mathrm{Re}\ln Y_r^k \end{array} \right), \mathbf{X} = \left(\begin{array}{c} \mathrm{Im}\ln Y_1^k \\ \mathrm{Im}\ln Y_2^k \\ \vdots \\ \mathrm{Im}\ln Y_r^k \end{array} \right)$$

Here $\mathbf{I}_r$ denotes the matrix identity of size r, where we denote

$$Y_i^k = \left(\begin{array}{c} \phi^*(0, -s_i; k) \\ \phi^*(s_i, 0; k) \\ \phi^*(s_i, -s_i; k) \end{array} \right) = \left(\begin{array}{c} \phi_1(s_i, k) \\ \phi_2(s_i, k) \\ \phi_3(s_i, k) \end{array} \right)$$

and the logarithm function is defined componentwise, i.e., we have

$$\mathrm{Re}\ln Y_i^k = \left(\begin{array}{c} \mathrm{Re}\ln \phi_1(s_i, k) \\ \mathrm{Re}\ln \phi_2(s_i, k) \\ \mathrm{Re}\ln \phi_3(s_i, k) \end{array} \right)$$

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and similarly for the imaginary part. Let us denote

$$\mathrm{E}Y_i^k = \left(\begin{array}{c} \mathrm{E}\phi_1(s_i,k) \\ \mathrm{E}\phi_2(s_i,k) \\ \mathrm{E}\phi_3(s_i,k) \end{array} \right) = \left(\begin{array}{c} \Phi_1(s_i,k) \\ \Phi_2(s_i,k) \\ \Phi_3(s_i,k) \end{array} \right)$$

Notice that $\Phi(u,v;k) = \mathrm{E}(\exp(i(uX_{t+k} + vX_t)))$, $u,v \in \mathbb{R}$. Using mean value theorem, we can expand the codifference function into

$$\left(\begin{array}{c} \mathrm{Re}\,\hat{\tau}^*(\mathbf{s},k) \\ \mathrm{Im}\,\hat{\tau}^*(\mathbf{s},k) \end{array} \right) = \lambda\left\{ \mathbf{L}_1^k + \bar{\mathbf{L}}_2^k \mathbf{Z}_N^k \right\} \tag{28}$$

where

$$\mathbf{L}_1^k = \left(\begin{array}{c} \mathrm{Re}\,\mathbf{L}_1^k \\ \mathrm{Im}\,\mathbf{L}_1^k \end{array} \right), \quad \mathbf{Z}_N^k = \left(\begin{array}{c} \mathrm{Re}\,\mathbf{Z}_N^k \\ \mathrm{Im}\,\mathbf{Z}_N^k \end{array} \right) = \left(\begin{array}{c} \mathrm{Re}\varphi_N^k - \mathrm{Re}\psi_N^k \\ \mathrm{Im}\varphi_N^k - \mathrm{Re}\psi_N^k \end{array} \right)$$

with

$$\mathrm{Re}\,\mathbf{L}_1^k = \left(\begin{array}{c} \mathrm{Re}\ln\mathrm{E}Y_1^k \\ \mathrm{Re}\ln\mathrm{E}Y_2^k \\ \vdots \\ \mathrm{Re}\ln\mathrm{E}Y_r^k \end{array} \right), \mathrm{Re}\varphi_N^k = \left(\begin{array}{c} \mathrm{Re}Y_1^k \\ \mathrm{Re}Y_2^k \\ \vdots \\ \mathrm{Re}Y_r^k \end{array} \right), \mathrm{Re}\psi_N^k = \left(\begin{array}{c} \mathrm{Re}\mathrm{E}Y_1^k \\ \mathrm{Re}\mathrm{E}Y_2^k \\ \vdots \\ \mathrm{Re}\mathrm{E}Y_r^k \end{array} \right)$$

and similarly for the imaginary parts, and where and $\bar{\mathbf{L}}_2^k = (\bar{d}_{ij}^k)_{i,j=1,\dots,6}$ denotes Jacobian of (26), which is evaluated at $\mathbf{c}$ ($\|\mathbf{c} - \psi_N^k\| < \|\varphi_N^k - \psi_N^k\|$). From the assumption C2, we obtain

$$\Phi_3(s_i,k) = \Phi(s_i,-s_i;k) = \exp\left(-\sum_{j=0}^{k-1}\sigma^\alpha |s_i c_j|^\alpha - \sum_{j=0}^{\infty}\sigma^\alpha |s_i(c_{j+k} - c_j)|^\alpha\right) \tag{29}$$

and $\Phi_1(s_i,k) = \Phi_2(s_i,k)$, i.e.,

$$\Phi(s_i,0;k) = \Phi(0,-s_i;k) = \exp\left(-\sum_{j=0}^{\infty}\sigma^\alpha |s_i c_j|^\alpha\right) \tag{30}$$

From identities (29)-(30) and further applying the assumption C1, we obtain that the elements of $\mathrm{Re}\,\psi_N^k$ are always strictly greater than 0. Therefore, with a probability convergent to 0, the elements of $\mathrm{Re}\,\varphi_N^k$ will be less than or equal to 0. Hence, without changing the limiting distribution of the estimator, we can restrict the definition of the real and the imaginary components of $\left(\begin{array}{c} \mathbf{Y} \\ \mathbf{X} \end{array} \right)$ in (26) only in the right half plane where the elements of $\mathrm{Re}(\varphi_N^k) > 0$, and equal to 0 in the other case. Thus, we can conclude that the Jacobian matrix $\bar{\mathbf{L}}_2^k$ is well defined here. By Theorem 2.1, $\bar{\mathbf{L}}_2^k$ will converge in probability to $\mathbf{L}_2^k$, where

$$\mathbf{L}_2^k = \nabla\mathbf{L}_1^k$$

Here ∇g denotes the Jacobian of g. From (29), (30), we have the following identities

$$\mathrm{Re}\,\Phi(s_i,-s_i;k) = \mathrm{E}\cos(s_i(X_{t+k} - X_t)) = \Phi(s_i,-s_i,k) \tag{31}$$

$$\mathrm{Re}\,\Phi(s_i,0;k) = \mathrm{E}\cos(s_i X_{t+k}) = \Phi(s_i,0,k) \tag{32}$$

$$\mathrm{Re}\,\Phi(0,-s_i;k) = \mathrm{E}\cos(-s_i X_t) = \Phi(0,-s_i,k) \tag{33}$$

and $\mathrm{Im}\,\Phi(s_i,-s_i;k) = \mathrm{E}\sin(s_i(X_{t+k}-X_t)) = 0$, $\mathrm{Im}\,\Phi(s_i,0;k) = \mathrm{E}\sin(s_i X_{t+k}) = 0$ and $\mathrm{Im}\,\Phi(0,-s_i;k) = \mathrm{E}\sin(-s_i X_t) = 0$. Using these identities, after some algebra we directly obtain

$$\mathbf{L}_2^k = \left(\begin{array}{cc} \mathbf{I}_r \mathbf{d}^k & \mathbf{0} \\ \mathbf{0} & \mathbf{I}_r \mathbf{d}^k \end{array} \right) \tag{34}$$

where $(\mathbf{d}^k)^T = [\mathbf{d}_1^k, \mathbf{d}_2^k, \ldots, \mathbf{d}_r^k]$, and the elements of $\mathbf{d}_i^k, i = 1, \ldots r$ are

$$d_i^k(1,1) = (\operatorname{Re} \Phi(0, -s_i; k))^{-1}$$
$$d_i^k(2,2) = (\operatorname{Re} \Phi(s_i, 0; k))^{-1}$$
$$d_i^k(3,3) = (\operatorname{Re} \Phi(s_i, -s_i; k))^{-1})$$

and equal to 0, otherwise. The asymptotic variance-covariance matrix is obtained from (28) as

$$\lim_{N \to \infty} N \ \operatorname{cov}\left(\left(\begin{array}{c} \operatorname{Re} \hat{\tau}^*(s, -s; p) \\ \operatorname{Im} \hat{\tau}^*(s, -s; p) \end{array} \right), \left(\begin{array}{c} \operatorname{Re} \hat{\tau}^*(s, -s; q) \\ \operatorname{Im} \hat{\tau}^*(s, -s; q) \end{array} \right)\right) = \lambda \mathbf{L}_2^p \mathbf{V}_{pq} \mathbf{L}_2^q \lambda^T \tag{35}$$

where

$$\mathbf{V}_{pq} = \left(\begin{array}{cc} \mathbf{V}_{pq}^{RR} & \mathbf{V}_{pq}^{RI} \\ \mathbf{V}_{pq}^{IR} & \mathbf{V}_{pq}^{II} \end{array} \right) = \lim_{N \to \infty} N \ \left(\begin{array}{cc} \operatorname{cov}(\operatorname{Re} \mathbf{Z}_N^p, \operatorname{Re} \mathbf{Z}_N^q) & \operatorname{cov}(\operatorname{Re} \mathbf{Z}_N^p, \operatorname{Im} \mathbf{Z}_N^q) \\ \operatorname{cov}(\operatorname{Im} \mathbf{Z}_N^p, \operatorname{Re} \mathbf{Z}_N^q) & \operatorname{cov}(\operatorname{Im} \mathbf{Z}_N^p, \operatorname{Im} \mathbf{Z}_N^q) \end{array} \right) \tag{36}$$

The matrix $\mathbf{V}_{pq}$ can be obtained by applying Theorem 1 and Remark 2.6. in Hesse (1990). Its elements can be derived in a similar way as obtaining variance-covariance matrix in Theorem 1 of Hesse (1990). This is possible, because it can be shown that all elements of $\mathbf{V}_{pq}$ (in the form of sum of the absolute components) are finite. Therefore, one can apply the property of the sample mean of ergodic processes (e.g., Theorem 7.1.1. in Brockwell and Davis, 1987). Notice that here in particular, we obtain all elements of $\mathbf{V}_{pq}$ with respect to $\operatorname{cov}(\operatorname{Re} \mathbf{Z}_N^p, \operatorname{Im} \mathbf{Z}_N^q)$ and $\operatorname{cov}(\operatorname{Im} \mathbf{Z}_N^p, \operatorname{Re} \mathbf{Z}_N^q)$ are zeros. The elements of $\mathbf{V}_{pq}$ with respect to $\operatorname{cov}(\operatorname{Re} \mathbf{Z}_N^p, \operatorname{Re} \mathbf{Z}_N^q)$ and $\operatorname{cov}(\operatorname{Im} \mathbf{Z}_N^p, \operatorname{Im} \mathbf{Z}_N^q)$ can be shown to be finite using identities (29)-(30) and applying a similar approach as obtaining eq. (21) and (23), and further applying Theorem A.1, or sometimes, eq.(2.7) in Kokoszka and Taqqu (1994) together with the similar steps as the proof of Theorem A.1. However, we omit details. ∎

Proposition B.2 *Let $X_t, t \in \mathbb{Z}$ be the moving average process of order m, $X_t = \sum_{j=0}^m c_j \epsilon_{t-j}$, satisfying conditions C1 and C2. Then for $h \in \{1, 2, \ldots\}$, $s \in \mathbb{R}, s \neq 0$*

$$\left[\left(\begin{array}{c} \operatorname{Re} \hat{\tau}^*(\mathbf{s}, 0) \\ \operatorname{Im} \hat{\tau}^*(\mathbf{s}, 0) \end{array} \right), \ldots, \left(\begin{array}{c} \operatorname{Re} \hat{\tau}^*(\mathbf{s}, h) \\ \operatorname{Im} \hat{\tau}^*(\mathbf{s}, h) \end{array} \right)\right] \text{ is } AN\left(\left[\left(\begin{array}{c} \tau(\mathbf{s}, 0) \\ 0 \end{array} \right), \ldots, \left(\begin{array}{c} \tau(\mathbf{s}, h) \\ 0 \end{array} \right)\right], N^{-1} M\right)$$

where $\mathbf{M}$ is the covariance matrix

$$\mathbf{M} = \left[\lambda \mathbf{L}_2^p \mathbf{V}_{pq} \mathbf{L}_2^q \lambda^T\right]_{p,q=0,\ldots,h}$$

and the matrices $\lambda, \mathbf{L}_2^k, k = p, q$ and $\mathbf{V}_{pq}$ are as given in Proposition B.1 above.

Proof. To show this relation, define vectors $\{\mathbf{Y}_t\}$ by

$$\mathbf{Y}_t^T = (\mathbf{Z}_t, \mathbf{Z}_{t+1}, \ldots, \mathbf{Z}_{t+h})$$

where

$$\mathbf{Z}_{t+k} = \left(\begin{array}{c} \mathbf{X}_1^k \\ \mathbf{X}_2^k \\ \vdots \\ \mathbf{X}_r^k \end{array} \right)$$

where for $j = 1, \ldots, r$

$$\mathbf{X}_j^k = \left(\begin{array}{c} \exp(-is_j X_t) \\ \exp(is_j X_{t+k}) \\ \exp(is_j (X_{t+k} - X_t)) \end{array} \right)$$

By definition, $\{\mathbf{Z}_{t+k}\}$ is $m+k$-dependent sequence and therefore $\{\mathbf{Y}_t\}$ is $m+h$-dependent sequence. Next define

$$\zeta_t^T = (\xi_t, \xi_{t+1}, \ldots, \xi_{t+h})$$

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where

$$\xi_{t+j} = \begin{pmatrix} \mathrm{Re}(\ln N^{-1} \sum_{t=1}^{N} Z_{t+j}) \\ \mathrm{Im}(\ln N^{-1} \sum_{t=1}^{N} Z_{t+j}) \end{pmatrix}$$

and

$$\mathrm{Re}(\ln N^{-1} \sum_{t=1}^{N} Z_{t+j}) = \begin{pmatrix} \mathrm{Re}(\ln N^{-1} \sum_{t=1}^{N} \mathbf{X}_1^j) \\ \mathrm{Re}(\ln N^{-1} \sum_{t=1}^{N} \mathbf{X}_2^j) \\ \vdots \\ \mathrm{Re}(\ln N^{-1} \sum_{t=1}^{N} \mathbf{X}_r^j)) \end{pmatrix}$$

where $l = 1, \ldots, r$

$$\mathrm{Re}(\ln N^{-1} \sum_{t=1}^{N} \mathbf{X}_l^j) = \begin{pmatrix} \mathrm{Re}(\ln N^{-1} \sum_{t=1}^{N} \exp(-is_l X_t)) \\ \mathrm{Re}(\ln N^{-1} \sum_{t=1}^{N} \exp(is_l X_{t+j})) \\ \mathrm{Re}(\ln N^{-1} \sum_{t=1}^{N} \exp(is_l (X_{t+j} - X_t))) \end{pmatrix}$$

(similarly for the imaginary part. Note that the summation and the principal value of $\ln(\cdot)$ are defined componentwise), then we have

$$\lambda \begin{pmatrix} \mathrm{Re} \ln (N^{-1} \sum_{t=1}^{N} \mathbf{Y}_t^T) \\ \mathrm{Im} \ln (N^{-1} \sum_{t=1}^{N} \mathbf{Y}_t^T) \end{pmatrix} = \lambda \zeta_t^T = \left[\begin{pmatrix} \mathrm{Re}\, \hat{\tau}^*(\mathbf{s}, 0) \\ \mathrm{Im}\, \hat{\tau}^*(\mathbf{s}, 0) \end{pmatrix}, \ldots, \begin{pmatrix} \mathrm{Re}\, \hat{\tau}^*(\mathbf{s}, h) \\ \mathrm{Im}\, \hat{\tau}^*(\mathbf{s}, h) \end{pmatrix} \right]$$

where λ is as given in (27). We therefore need to show that when $N \to \infty$

$$\mathbf{a}^T (\lambda[\xi_t, \xi_{t+1}, \ldots, \xi_{t+h}])^T \text{ is } AN \left(\mathbf{a}^T \left(\begin{pmatrix} \mathrm{Re}\, \tau(0) \\ 0 \end{pmatrix}, \ldots, \begin{pmatrix} \mathrm{Re}\, \tau(h) \\ 0 \end{pmatrix} \right)^T, N^{-1} \mathbf{a}^T \mathbf{M} \mathbf{a} \right) \quad (37)$$

for all vectors $\mathbf{a} = (a_0, \ldots, a_h)^T \in \mathbb{R}^{h+1}$ such that $\mathbf{a}^T \mathbf{M} \mathbf{a} > 0$. For any such $\mathbf{a}$, the sequence $\{\mathbf{a}^T (\lambda \zeta_t^T)^T\}$ is $(m+h)$-dependent and since by Proposition B.1

$$\lim_{N \to \infty} N \; \mathrm{var}(\mathbf{a}^T (\lambda[\xi_t, \xi_{t+1}, \ldots, \xi_{t+h}])^T) = \mathbf{a}^T \mathbf{M} \mathbf{a} > 0$$

where $\mathbf{M}$ is the covariance matrix

$$\mathbf{M} = \left[\lambda \mathbf{L}_2^p \mathbf{V}_{pq} \mathbf{L}_2^q \lambda^T \right]_{p,q=0,\ldots,h}$$

and the vectors $\lambda, \mathbf{L}_2^p, \mathbf{L}_2^q$, matrix $\mathbf{V}_{pq}$ are as given in Proposition B.1 above. We can conclude that $\{\mathbf{a}^T (\lambda \zeta_t^T)^T\}$ satisfies the conditions of central limit theorem for m-dependent processes (e.g Brockwell and Davis, 1987, Theorem 6.4.2), and therefore by this theorem, for $N \to \infty$, we obtain the required result (37). The relation $\mathrm{Im}\, \tau(\mathbf{s}, j) = 0, j = 0, 1, \ldots, h$ can be obtained directly from identities (29)-(30). ∎

Proposition B.3 *Proposition B.2 remains true for $X_t, t \in \mathbb{Z}$ being a stationary linear process (1), satisfying conditions C1 and C2.*

Proof. For the proof, we will apply the result of Proposition B.2 to the truncated sequence $X_{tm} = \sum_{j=0}^{m} c_j \epsilon_{t-j}$ and then derive the result for X_t by letting $m \to \infty$. For $0 \le p \le h$, we define

$$\hat{\tau}_m^*(s, -s; p) = -\ln \phi_m^*(s, -s; p) + \ln \phi_m^*(s, 0; p) + \ln \phi_m^*(0, -s; p) \quad (38)$$

where $\phi_m^*(u, v; p) = N^{-1} \sum_{t=1}^{N} \exp(i(u X_{(t+p)m} + v X_{tm}))$. Then by Proposition B.2

$$N^{1/2} \left[\begin{pmatrix} \mathrm{Re}\, \hat{\tau}_m^*(\mathbf{s}, 0) - \mathrm{Re}\, \tau_m(\mathbf{s}, 0) \\ \mathrm{Im}\, \hat{\tau}_m^*(\mathbf{s}, 0) - \mathrm{Im}\, \tau_m(\mathbf{s}, 0) \end{pmatrix}, \ldots, \begin{pmatrix} \mathrm{Re}\, \hat{\tau}_m^*(\mathbf{s}, h) - \mathrm{Re}\, \tau_m(\mathbf{s}, h) \\ \mathrm{Im}\, \hat{\tau}_m^*(\mathbf{s}, h) - \mathrm{Im}\, \tau_m(\mathbf{s}, h) \end{pmatrix} \right] \Rightarrow Y_m$$

where $Y_m \sim N(0, M_m)$. Here $\mathbf{M}_m$ is the covariance matrix

$$\mathbf{M}_m = \left(\begin{array}{cc} \text{cov}(\text{Re}\,\hat{\tau}_m^*(\mathbf{s}, p), \text{Re}\,\hat{\tau}_m^*(\mathbf{s}, q)) & \text{cov}(\text{Re}\,\hat{\tau}_m^*(\mathbf{s}, p), \text{Im}\,\hat{\tau}_m^*(\mathbf{s}, q)) \\ \text{cov}(\text{Im}\,\hat{\tau}_m^*(\mathbf{s}, p), \text{Re}\,\hat{\tau}_m^*(\mathbf{s}, q)) & \text{cov}(\text{Im}\,\hat{\tau}_m^*(\mathbf{s}, p), \text{Im}\,\hat{\tau}_m^*(\mathbf{s}, q)) \end{array} \right)_{p,q=0,\ldots,h}$$

$$= \left[\lambda \mathbf{L}_2^p(m) \mathbf{V}_{pq}^m \mathbf{L}_2^q(m) \lambda^T \right]_{p,q=0,\ldots,h}$$

where λ is defined as (27) and the Jacobian matrix $\mathbf{L}_2^k(m)$ and matrix $\mathbf{V}_{pq}^m$ are defined for X_{tm} as in (34) and (36), respectively. Now, as $m \to \infty$,

$$\mathbf{M}_m \to \mathbf{M}$$

where $\mathbf{M}$ is defined like $\mathbf{M}_m$ by replacing X_{tm} by X_t. Hence

$$Y_m \Rightarrow Y \text{ where } Y \sim N(0, \mathbf{M})$$

The proof now can be completed by applying Proposition 6.3.9. in Brockwell and Davis (1987) provided we can show that

$$\lim_{m \to \infty} \limsup_{N \to \infty} \text{P}(N^{1/2} \left| \text{Re}\,\hat{\tau}_m^*(\mathbf{s}, p) - \text{Re}\,\tau_m(\mathbf{s}, p) - \text{Re}\,\hat{\tau}^*(\mathbf{s}, p) + \text{Re}\,\tau(\mathbf{s}, p) \right| > \epsilon) = 0 \qquad (39)$$

for $p = 0, 1, \ldots, h$ (and similarly for the imaginary part). The probability in (39) is bounded by

$$\epsilon^{-2} N \,\text{var}(\text{Re}\,\hat{\tau}_m^*(\mathbf{s}, p) - \text{Re}\,\hat{\tau}^*(\mathbf{s}, p))$$

$$= \epsilon^{-2} \left[N \,\text{var}(\text{Re}\,\hat{\tau}_m^*(\mathbf{s}, p)) + N \,\text{var}(\text{Re}\,\hat{\tau}^*(\mathbf{s}, p)) - 2N\text{cov}(\text{Re}\,\hat{\tau}_m^*(\mathbf{s}, p), \text{Re}\,\hat{\tau}^*(\mathbf{s}, p)) \right]$$

From the calculation of Proposition B.1 and further noting that Theorem 1 and Remark 2.6. in Hesse (1990) can be applied for the finite moving average process by setting some of the coefficients c_j's to be zero, we obtain

$$\lim_{m \to \infty} \lim_{N \to \infty} N \,\text{var}(\text{Re}\,\hat{\tau}_m^*(\mathbf{s}, p)) = \lim_{N \to \infty} N \,\text{var}(\text{Re}\,\hat{\tau}^*(\mathbf{s}, p)) = m_{pp}^{RR}$$

where m_{pq}^{RR} denotes the covariance between the real elements in (p, q)- block of covariance matrix $\mathbf{M}$. Moreover, using the same steps to that given in the proof of Proposition B.1, it can be shown that

$$\lim_{m \to \infty} \lim_{N \to \infty} N\text{cov}(\text{Re}\,\hat{\tau}_m^*(\mathbf{s}, p), \text{Re}\,\hat{\tau}^*(\mathbf{s}, p) = m_{pp}^{RR}$$

Thus

$$\lim_{m \to \infty} \limsup_{N \to \infty} \epsilon^{-2} N \,\text{var}(\text{Re}\,\hat{\tau}_m^*(\mathbf{s}, p) - \text{Re}\,\hat{\tau}^*(\mathbf{s}, p)) = 0$$

Similar results can be obtained for the imaginary part. This established (39). ∎

Theorem B.4 *Let $X_t, t \in \mathbb{Z}$ be the stationary linear process (1), satisfying conditions C1 and C2. Then for $h \in \{1, 2, \ldots\}$, $s \in \mathbb{R}, s \neq 0$*

$$\left[\left(\begin{array}{c} \text{Re}\,\hat{\tau}(\mathbf{s}, 0) \\ \text{Im}\,\hat{\tau}(\mathbf{s}, 0) \end{array} \right), \ldots, \left(\begin{array}{c} \text{Re}\,\hat{\tau}(\mathbf{s}, h) \\ \text{Im}\,\hat{\tau}(\mathbf{s}, h) \end{array} \right) \right] \text{ is } AN \left(\left[\left(\begin{array}{c} \tau(\mathbf{s}, 0) \\ 0 \end{array} \right), \ldots, \left(\begin{array}{c} \tau(\mathbf{s}, h) \\ 0 \end{array} \right) \right], N^{-1}M \right)$$

where $\mathbf{M}$ is as given in Proposition B.2 above.

Proof. To show the convergence of the estimator $\text{Re}\,\hat{\tau}(\mathbf{s}, j)$ and $\text{Im}\,\hat{\tau}(\mathbf{s}, j)$ to the same limit as $\text{Re}\,\hat{\tau}^*(\mathbf{s}, j)$ and $\text{Im}\,\hat{\tau}^*(\mathbf{s}, j)$, respectively, with $0 \leq j \leq h$, it suffices to show that as $N \to \infty$

$$N^{1/2} \left\{ \lambda_2 \left(\begin{array}{c} \text{Re}\,\phi^*(s_k, -s_k; j) \\ \text{Re}\,\phi^*(s_k, 0; j) \\ \text{Re}\,\phi^*(0, -s_k; j) \end{array} \right) - \lambda_2 \left(\begin{array}{c} \text{Re}\,\phi(s_k, -s_k; j) \\ \text{Re}\,\phi(s_k, 0; j) \\ \text{Re}\,\phi(0, -s_k; j) \end{array} \right) \right\} = o_p(1)$$

(and similarly for the imaginary part), where $\phi^*(u,v;j) = N^{-1}\sum_{t=1}^{N}\exp(i(uX_{t+j}+vX_t))$, $\phi(u,v;j)=$ $(N-j)^{-1}\sum_{t=1}^{N-j}\exp(i(uX_{t+j}+vX_t))$ and $\lambda_2 = \begin{bmatrix} -1 & 1 & 1 \end{bmatrix}$. The required result then follows from Slutzky's theorem (e.g., Theorem 5.1.1. in Lehmann, 1999).

Simple algebra gives, for $0 \le j \le h$,

$$
N^{1/2}\mathrm{E}\left| \lambda_2 \begin{pmatrix} \mathrm{Re}\,\phi^*(s_k,-s_k;j) \\ \mathrm{Re}\,\phi^*(s_k,0;j) \\ \mathrm{Re}\,\phi^*(0,-s_k;j) \end{pmatrix} - \lambda_2 \begin{pmatrix} \mathrm{Re}\,\phi(s_k,-s_k;j) \\ \mathrm{Re}\,\phi(s_k,0;j) \\ \mathrm{Re}\,\phi(0,-s_k;j) \end{pmatrix} \right|
$$

$$
= N^{1/2}\mathrm{E}\left| \lambda_2 \begin{pmatrix} \frac{j}{(N-j)}\frac{1}{N}\sum_{t=1}^{N}\cos(is_k(X_{t+j}-X_t)) - \frac{1}{N-j}\sum_{t=N-j+1}^{N}\cos(is(X_{t+j}-X_t)) \\ \frac{j}{(N-j)}\frac{1}{N}\sum_{t=1}^{N}\cos(is_k X_{t+j}) - \frac{1}{N-j}\sum_{t=N-j+1}^{N}\cos(is_k X_{t+j}) \\ \frac{j}{(N-j)}\frac{1}{N}\sum_{t=1}^{N}\cos(-is_k X_t) - \frac{1}{N-j}\sum_{t=N-j+1}^{N}\cos(-is_k X_t) \end{pmatrix} \right|
$$

$$
\le 6j(N-j)^{-1/2}\left(\tfrac{N}{N-j}\right)^{1/2}
$$

The required result is obtained from $3j(N-j)^{-1/2} \to 0$ and $N/(N-j) \to 1$ as $N \to \infty$. Using the same arguments, similar results can be obtained for the imaginary part. The conclusion of the theorem then follows from Proposition B.3 above. ∎

Proof of Theorem 2.2. Let $\mathbf{g}(\cdot)$ be the function from $\mathbb{R}^{2r\times(h+1)}$ to $\mathbb{R}^{2r\times h}$ defined by

$$
\mathbf{g}\left(\left[\begin{pmatrix} \hat\tau(\mathbf{s},0) \\ 0 \end{pmatrix}, \begin{pmatrix} \mathrm{Re}\,\hat\tau(\mathbf{s},1) \\ \mathrm{Im}\,\hat\tau(\mathbf{s},1) \end{pmatrix}, \dots, \begin{pmatrix} \mathrm{Re}\,\hat\tau(\mathbf{s},h) \\ \mathrm{Im}\,\hat\tau(\mathbf{s},h) \end{pmatrix} \right]^T \right)
$$

$$
= \left[\begin{pmatrix} \mathrm{Re}\,\hat I(\mathbf{s},1) \\ \mathrm{Im}\,\hat I(\mathbf{s},1) \end{pmatrix}, \begin{pmatrix} \mathrm{Re}\,\hat I(\mathbf{s},2) \\ \mathrm{Im}\,\hat I(\mathbf{s},2) \end{pmatrix}, \dots, \begin{pmatrix} \mathrm{Re}\,\hat I(\mathbf{s},h) \\ \mathrm{Im}\,\hat I(\mathbf{s},h) \end{pmatrix} \right]^T
$$

where for $0 < j \le h$ and $\hat\tau(0) \ne 0$, we have $\mathrm{Re}\,\hat I(s_i,-s_i;j) = \frac{\mathrm{Re}\,\hat\tau(s_i,-s_i;j)}{\hat\tau(s_i,-s_i;0)}$ and $\mathrm{Im}\,\hat I(s_i,-s_i;j) = \frac{\mathrm{Im}\,\hat\tau(s_i,-s_i;j)}{\hat\tau(s_i,-s_i;0)}$, for $i = 1,\dots,r$. By applying delta method and Theorem B.4 above, we can show that

$$
\left[\begin{pmatrix} \mathrm{Re}\,\hat I(\mathbf{s},1) \\ \mathrm{Im}\,\hat I(\mathbf{s},1) \end{pmatrix}, \begin{pmatrix} \mathrm{Re}\,\hat I(\mathbf{s},2) \\ \mathrm{Im}\,\hat I(\mathbf{s},2) \end{pmatrix}, \dots, \begin{pmatrix} \mathrm{Re}\,\hat I(\mathbf{s},h) \\ \mathrm{Im}\,\hat I(\mathbf{s},h) \end{pmatrix} \right]^T
$$

is asymptotically normal distributed with mean

$$
\mathbf{g}\left(\left[\begin{pmatrix} \tau(\mathbf{s},0) \\ 0 \end{pmatrix}, \begin{pmatrix} \tau(\mathbf{s},1) \\ 0 \end{pmatrix}, \dots, \begin{pmatrix} \tau(\mathbf{s},h) \\ 0 \end{pmatrix} \right]^T \right)
$$

$$
= \left[\begin{pmatrix} I(1) \\ 0 \end{pmatrix}, \begin{pmatrix} I(2) \\ 0 \end{pmatrix}, \dots, \begin{pmatrix} I(h) \\ 0 \end{pmatrix} \right]^T
$$

and variance $N^{-1}\mathbf{D}\mathbf{M}\mathbf{D}^{\mathbf{T}}$. Here the matrix $\mathbf{M}$ is as given in Proposition B.2 and $\mathbf{D}$ is the Jacobian matrix of $\mathbf{g}(\cdot)$. To obtain the elements of matrix $\mathbf{D}$, we proceed as follows. First, note that the codifference function at lag 0 is a real-valued function. Therefore, for $0 \le j \le h$, and $\tau(0) \ne 0$, we obtain $\mathrm{Re}\,I(j) = \frac{\mathrm{Re}\,\tau(j)}{\tau(0)} = I(j)$ and $\mathrm{Im}\,I(j) = \frac{\mathrm{Im}\,\tau(j)}{\tau(0)} = 0$. It is straightforward to obtain the Jacobian matrix $\mathbf{D}$ as

$$
\mathbf{D} = \begin{bmatrix} D_{11} & D_{12} & 0 & \dots & 0 \\ D_{21} & 0 & D_{23} & & 0 \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ D_{h1} & 0 & 0 & \dots & D_{h(h+1)} \end{bmatrix} \tag{40}
$$

where

$$
D_{l1} = \begin{bmatrix} D_{l1}^{11} & \mathbf{0}_r \\ \mathbf{0}_r & \mathbf{0}_r \end{bmatrix} \tag{41}
$$

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and

$$D_{l(l+1)} = \begin{bmatrix} D_{l(l+1)}^{11} & \mathbf{0}_r \\ \mathbf{0}_r & D_{l(l+1)}^{11} \end{bmatrix} \qquad (42)$$

for $l = 1, ..h$, where

$$D_{l1}^{11} = \mathbf{I}_r \left[\frac{-I(l)}{\tau(s_1, -s_1; 0)}, \frac{-I(l)}{\tau(s_2, -s_2; 0)}, \dots, \frac{-I(l)}{\tau(s_r, -s_r; 0)} \right]^T$$

and

$$D_{l(l+1)}^{11} = \mathbf{I}_r \left[\frac{1}{\tau(s_1, -s_1; 0)}, \frac{1}{\tau(s_2, -s_2; 0)}, \dots, \frac{1}{\tau(s_r, -s_r; 0)} \right]^T$$

Here $\mathbf{I}_r$ denotes the matrix identity of size r. Let's denote w_{ij}, for $i, j = 1, \dots, h$, the (i, j)-th block element of $\mathbf{DMD}^T$ and m_{ij}, for $i, j = 0, 1, \dots, h$, the (i, j)-th block element of $\mathbf{M}$. We find that

$$
\begin{aligned}
w_{ij} &= \begin{bmatrix} \mathrm{cov}(\mathrm{Re}\,\hat{I}(\mathbf{s}, i), \mathrm{Re}\,\hat{I}(\mathbf{s}, j)) & \mathrm{cov}(\mathrm{Re}\,\hat{I}(\mathbf{s}, i), \mathrm{Im}\,\hat{I}(\mathbf{s}, j)) \\ \mathrm{cov}(\mathrm{Im}\,\hat{I}(\mathbf{s}, i), \mathrm{Re}\,\hat{I}(\mathbf{s}, j)) & \mathrm{cov}(\mathrm{Im}\,\hat{I}(\mathbf{s}, i), \mathrm{Im}\,\hat{I}(\mathbf{s}, j)) \end{bmatrix} \\
&= D_{i1} m_{00} D_{j1} + D_{i(i+1)} m_{i0} D_{j1} + D_{i1} m_{0j} D_{j(j+1)} + D_{i(i+1)} m_{ij} D_{j(j+1)} \\
&= \begin{bmatrix} D_{i1}^{11} m_{00}^{RR} D_{j1}^{11} + D_{i(i+1)}^{11} m_{i0}^{RR} D_{j1}^{11} + D_{i1}^{11} m_{0j}^{RR} D_{j(j+1)}^{11} + D_{i(i+1)}^{11} m_{ij}^{RR} D_{j(j+1)}^{11} & \mathbf{0}_r \\ \mathbf{0}_r & D_{i(i+1)}^{11} m_{ij}^{II} D_{j(j+1)}^{11} \end{bmatrix}
\end{aligned}
$$
$$(43)$$

Here m_{ij}^{RR} and m_{ij}^{II} denote the partitions of m_{ij} which correspond to the real and the imaginary components, respectively. ∎

C   Proof of Corollary 2.3

Proof of Corollary 2.3. As MA(0) is a special case of the linear process (1), by applying Theorem 2.2, one can conclude the asymptotic normality of

$$\left[\begin{pmatrix} \mathrm{Re}\,\hat{I}(\mathbf{s}, 1) \\ \mathrm{Im}\,\hat{I}(\mathbf{s}, 1) \end{pmatrix}, \begin{pmatrix} \mathrm{Re}\,\hat{I}(\mathbf{s}, 2) \\ \mathrm{Im}\,\hat{I}(\mathbf{s}, 2) \end{pmatrix}, \dots, \begin{pmatrix} \mathrm{Re}\,\hat{I}(\mathbf{s}, h) \\ \mathrm{Im}\,\hat{I}(\mathbf{s}, h) \end{pmatrix} \right]^T$$

for $h \in \{1, 2, \dots\}$. The true codifference function of i.i.d. process X_t is

$$
\begin{aligned}
\tau(s, -s; k) &= -\ln \mathrm{E} \exp(is(X_{t+k} - X_t)) + \ln \mathrm{E} \exp(isX_{t+k}) + \ln \mathrm{E} \exp(-isX_t) \\
&= \begin{cases} -2\sigma^\alpha |s|^\alpha & \text{for } k = 0 \\ 0 & \text{for } k > 0 \end{cases}
\end{aligned}
$$

which enables us to conclude that the real and the imaginary parts of $I(k) = 0$ whenever $k > 0$.

From (43), we obtain that $w_{kk}, k > 0$ is reduced to

$$w_{kk} = D_{k(k+1)} m_{kk} D_{k(k+1)} \qquad (44)$$

where matrix $D_{k(k+1)}$ is as given in (42), with

$$D_{l(l+1)}^{11} = \mathbf{I}_r \left[\frac{1}{-2\sigma^\alpha |s_1|^\alpha}, \frac{1}{-2\sigma^\alpha |s_2|^\alpha}, \dots, \frac{1}{-2\sigma^\alpha |s_r|^\alpha} \right]^T$$

and where

$$m_{kk} = \lambda \mathbf{L}_2^k \mathbf{V}_{kk} \mathbf{L}_2^k \lambda^T \qquad (45)$$

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with λ is as given as (27), and the elements of the matrix $\mathbf{L}_2^k$ and the covariance matrix $\mathbf{V}_{kk}$ will be given below. Let us denote

$$\mathbf{V}_{kk}^{RR}(i,j) = [\text{cov}(\text{Re}\,\phi_p(s_i,k), \text{Re}\,\phi_q(s_j,k)))]_{p,q=1,2,3}$$

and

$$\mathbf{V}_{kk}^{II}(i,j) = [\text{cov}(\text{Im}\,\phi_p(s_i,k), \text{Im}\,\phi_q(s_j,k)))]_{p,q=1,2,3}$$

as the (i,j)-th block elements of $\mathbf{V}_{kk}^{RR}$ and $\mathbf{V}_{kk}^{II}$, respectively. Using identities (31)-(33) (and the identities for imaginary part afterwards) in p.12, we can obtain their components as follows

$$\text{cov}(\text{Re}(\phi_1(s_i,p)), \text{Re}(\phi_1(s_j,q))) = \text{cov}(\cos(-s_i X_t), \cos(-s_j X_t))$$
$$= \tfrac{1}{2}\{e^{-\sigma^\alpha|s_i+s_j|^\alpha} + e^{-\sigma^\alpha|s_i-s_j|^\alpha}\} - e^{-\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha)}$$

$$\text{cov}(\text{Re}(\phi_1(s_i,p)), \text{Re}(\phi_2(s_j,q))) = \text{cov}(\text{Re}(\phi_2(s_i,p)), \text{Re}(\phi_1(s_j,q))) = \text{cov}(\text{Re}(\phi_1(s_i,p)), \text{Re}(\phi_1(s_j,q)))$$

$$\text{cov}(\text{Re}(\phi_2(s_i,p)), \text{Re}(\phi_2(s_j,q))) = \text{cov}(\text{Re}(\phi_1(s_i,p)), \text{Re}(\phi_1(s_j,q)))$$

$$\text{cov}(\text{Re}(\phi_1(s_i,p)), \text{Re}(\phi_3(s_j,q)))$$
$$= \text{cov}(\cos(-s_i X_t), \cos(s_j(X_{t+q} - X_t))) + \text{cov}(\cos(-s_i X_{t+q}), \cos(s_j(X_{t+q} - X_t)))$$
$$= e^{-\sigma^\alpha(|s_j|^\alpha+|s_i-s_j|^\alpha)} + e^{-\sigma^\alpha(|s_j|^\alpha+|s_i+s_j|^\alpha)} - 2e^{-\sigma^\alpha(|s_i|^\alpha+|2s_j|^\alpha)}$$

$$\text{cov}(\text{Re}(\phi_3(s_i,p)), \text{Re}(\phi_1(s_j,q)))$$
$$= \text{cov}(\cos(-s_j X_t), \cos(s_i(X_{t+p} - X_t))) + \text{cov}(\cos(-s_j X_{t+p}), \cos(s_i(X_{t+p} - X_t)))$$
$$= e^{-\sigma^\alpha(|s_i|^\alpha+|s_i-s_j|^\alpha)} + e^{-\sigma^\alpha(|s_i|^\alpha+|s_i+s_j|^\alpha)} - 2e^{-\sigma^\alpha(|s_j|^\alpha+|2s_i|^\alpha)}$$

$$\text{cov}(\text{Re}(\phi_2(s_i,p)), \text{Re}(\phi_3(s_j,q)))$$
$$= \text{cov}(\cos(s_i X_{t+q}), \cos(s_j(X_{t+q} - X_t))) + \text{cov}(\cos(s_i X_{t+p}), \cos(s_j(X_{t+p+q} - X_{t+p})))$$
$$= e^{-\sigma^\alpha(|s_j|^\alpha+|s_i-s_j|^\alpha)} + e^{-\sigma^\alpha(|s_j|^\alpha+|s_i+s_j|^\alpha)} - 2e^{-\sigma^\alpha(|s_i|^\alpha+|2s_j|^\alpha)}$$

$$\text{cov}(\text{Re}(\phi_3(s_i,k)), \text{Re}(\phi_2(s_j,k)))$$
$$= \text{cov}(\cos(s_j X_{t+k}), \cos(s_i(X_{t+k} - X_t))) + \text{cov}(\cos(s_j X_{t+k}), \cos(s_i(X_{t+2k} - X_{t+k})))$$
$$= e^{-\sigma^\alpha(|s_i|^\alpha+|s_i-s_j|^\alpha)} + e^{-\sigma^\alpha(|s_i|^\alpha+|s_i+s_j|^\alpha)} - 2e^{-\sigma^\alpha(|s_j|^\alpha+|2s_i|^\alpha)}$$

$$\text{cov}(\text{Re}(\phi_3(s_i,p)), \text{Re}(\phi_3(s_j,q)))$$
$$= \text{cov}(\cos(s_i(X_{t+p} - X_t)), \cos(s_j(X_{t+q} - X_t))) + \text{cov}(\cos(s_i(X_{t+p+q} - X_{t+q})), \cos(s_j(X_{t+q} - X_t)))$$
$$\quad + \text{cov}(\cos(s_i(X_{t+p} - X_t)), \cos(s_j(X_{t++p+q} - X_{t+p}))) + c_{\text{Re}}^{pq}$$

where

$$c_{\text{Re}}^{pq} = \begin{cases} 0 & \text{if } p = q \\ \text{cov}(\cos(s_i(X_{t+q} - X_{t+q-p})), \cos(s_j(X_{t+q} - X_t))) & \text{if } q > p \\ \text{cov}(\cos(s_i(X_{t+p} - X_t)), \cos(s_j(X_{t+p} - X_{t+p-q}))) & \text{if } p > q \end{cases}$$

yielding for $p = q$

$$\text{cov}(\text{Re}(\phi_3(s_i,p)), \text{Re}(\phi_3(s_j,q))) = \tfrac{1}{2}e^{-2\sigma^\alpha|s_i+s_j|^\alpha} + \tfrac{1}{2}e^{-2\sigma^\alpha|s_i-s_j|^\alpha} - 3e^{-\sigma^\alpha(2|s_i|^\alpha+|2s_j|^\alpha)}$$
$$+ e^{-\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha+|s_i-s_j|^\alpha)} + e^{-\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha+|s_i+s_j|^\alpha)}$$

and for $p \neq q$

$$\mathrm{cov}(\mathrm{Re}(\phi_3(s_i,p)),\mathrm{Re}(\phi_3(s_j,q))) = 2e^{-\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha+|s_i-s_j|^\alpha)}$$
$$+ 2e^{-\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha+|s_i+s_j|^\alpha)} - 4e^{-\sigma^\alpha(2|s_i|^\alpha+|2s_j|^\alpha)}$$

$$\mathrm{cov}(\mathrm{Im}(\phi_1(s_i,p)),\mathrm{Im}(\phi_1(s_j,q))) = \mathrm{cov}(\sin(-s_iX_t),\sin(-s_jX_t)) = \tfrac{1}{2}\{e^{-\sigma^\alpha|s_i-s_j|^\alpha} - e^{-\sigma^\alpha|s_i+s_j|^\alpha}\}$$

$$\mathrm{cov}(\mathrm{Im}(\phi_1(s_i,p)),\mathrm{Im}(\phi_2(s_j,q))) = \mathrm{cov}(\mathrm{Im}(\phi_2(s_i,p)),\mathrm{Im}(\phi_1(s_j,q))) \quad = -\mathrm{cov}(\mathrm{Im}(\phi_1(s_i,p)),\mathrm{Im}(\phi_1(s_j,q)))$$

$$\mathrm{cov}(\mathrm{Im}(\phi_2(s_i,p)),\mathrm{Im}(\phi_2(s_j,q))) = \mathrm{cov}(\mathrm{Im}(\phi_1(s_i,p)),\mathrm{Im}(\phi_1(s_j,q)))$$

$$\mathrm{cov}(\mathrm{Im}(\phi_3(s_i,p)),\mathrm{Im}(\phi_3(s_j,q)))$$
$$= \mathrm{cov}(\sin(s_i(X_{t+p}-X_t)),\sin(s_j(X_{t+q}-X_t))) + \mathrm{cov}(\sin(s_i(X_{t+p+q}-X_{t+q})),\sin(s_j(X_{t+q}-X_t)))$$
$$+ \mathrm{cov}(\sin(s_i(X_{t+p+q}-X_{t+p})),\sin(s_j(X_{t+p}-X_t))) + c_{\mathrm{Im}}^{pq}$$

where

$$c_{\mathrm{Im}}^{pq} = \begin{cases} 0 & \text{if } p=q \\ \mathrm{cov}(\sin(s_i(X_{t+q}-X_{t+q-p})),\sin(s_j(X_{t+q}-X_t))) & \text{if } q>p \\ \mathrm{cov}(\sin(s_i(X_{t+p}-X_t)),\sin(s_j(X_{t+p}-X_{t+p-q}))) & \text{if } p>q \end{cases}$$

yielding for $p=q$

$$\mathrm{cov}(\mathrm{Im}(\phi_3(s_i,k)),\mathrm{Im}(\phi_3(s_j,k)))$$
$$= \tfrac{1}{2}e^{-2\sigma^\alpha|s_i-s_j|^\alpha} - \tfrac{1}{2}e^{-2\sigma^\alpha|s_i+s_j|^\alpha} + e^{-\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha+|s_i+s_j|^\alpha)} - e^{-\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha+|s_i-s_j|^\alpha)}$$

and $\mathrm{cov}(\mathrm{Im}(\phi_3(s_i,k)),\mathrm{Im}(\phi_3(s_j,k))) = 0$ for $p \neq q$. The other elements are all zeros. The elements of $\mathbf{L}_2^k$ are as given in (34), where the elements of $\mathbf{d}_i^k, i=1,\ldots r$ are

$$d_i^k(1,1) = (\mathrm{Re}\,\Phi(0,-s_i;k))^{-1} = e^{\sigma^\alpha|s_i|^\alpha}$$
$$d_i^k(2,2) = (\mathrm{Re}\,\Phi(s_i,0;k))^{-1} = e^{\sigma^\alpha|s_i|^\alpha}$$
$$d_i^k(3,3) = (\mathrm{Re}\,\Phi(s_i,-s_i;k))^{-1} = e^{2\sigma^\alpha|s_i|^\alpha})$$

As from (45) we obtain

$$m_{kk}^{RR} = \mathrm{cov}(\mathrm{Re}\,\hat\tau(\mathbf{s},k),\mathrm{Re}\,\hat\tau(\mathbf{s},k)) = (\mathbf{I}_r \otimes \lambda_1)d^k V_{kk}^{RR} d^k(\mathbf{I}_r \otimes \lambda_1^T)$$

and

$$m_{kk}^{II} = \mathrm{cov}(\mathrm{Im}\,\hat\tau(\mathbf{s},k),\mathrm{Im}\,\hat\tau(\mathbf{s},k)) = (\mathbf{I}_r \otimes \lambda_1)d^k V_{kk}^{RR} d^k(\mathbf{I}_r \otimes \lambda_1^T)$$

then the (i,j)-th element of m_{kk}^{RR} and m_{kk}^{II} is obtained from

$$m_{kk}^{RR}(i,j) = \lambda_1 d_i^k V_{kk}^{RR}(i,j)d_j^k \lambda_1^T$$

and

$$m_{kk}^{II}(i,j) = \lambda_1 d_i^k V_{kk}^{II}(i,j)d_j^k \lambda_1^T$$

which therefore after a simple algebra, we obtain

$$m_{kk}^{RR}(i,j) = e^{\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha-|s_i-s_j|^\alpha)}\left\{\tfrac{1}{2}e^{\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha-|s_i-s_j|^\alpha)} - 1\right\}$$
$$+ e^{\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha-|s_i+s_j|^\alpha)}\left\{\tfrac{1}{2}e^{\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha-|s_i+s_j|^\alpha)} - 1\right\} + 1$$

$$m_{kk}^{II}(i,j) = e^{\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha-|s_i-s_j|^\alpha)}\left\{\tfrac{1}{2}e^{\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha-|s_i-s_j|^\alpha)} - 1\right\}$$
$$+ e^{\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha-|s_i+s_j|^\alpha)}\left\{1 - \tfrac{1}{2}e^{\sigma^\alpha(|s_i|^\alpha+|s_j|^\alpha-|s_i+s_j|^\alpha)}\right\}$$

The required result follows directly from (44). ∎

References

Adler, R. J., R. E. Feldman and C. M. Gallagher (1998). Analyzing stable time series. In: *A Practical Guide to Heavy Tails : Statistical Techniques and Applications* (R. J. Adler, R. E. Feldman and M. S. Taqqu, Eds.). pp. 133–158. Birkhäuser. Boston.

Bhansali, R. J. (1983). Estimation for the order of MA model from autoregressive and windows estimates of the inverse correlation function. *Journal of Time Series Analysis* **4**, 137–162.

Brockwell, P. J. and R. A. Davis (1987). *Time Series: Theory and Methods*. Springer-Verlag. New York.

Chambers, J. M., C. L. Mallows and B. W. Stuck (1976). A method for simulating stable random variables. *Jounal of the American Statistical Association* **71**, 340–344.

Embrechts, P., C. Klüppelberg and T. Mikosch (1997). *Modelling Extremal Events for Insurance and Finance*. Springer Verlag. Berlin.

Hesse, C. H. (1990). Rates of convergence for the empirical distribution function and the empirical characteristic function of a broad class of linear process. *Journal of Multivariate Analysis* **35**, 186–202.

Hong, Y. (1999). Hypothesis testing in time series via the empirical characteristic function: A generalized spectral density approach. *Journal of the American Statistical Association* **94**(448), 1201–1214.

Janicki, A. and A. Weron (1994). *Simulation and Chaotic Behavior of α-Stable Stochastic Processes*. Marcel Dekker. New York.

Kogon, S. M. and D. B. Williams (1998). Characteristic function based estimation of stable distribution parameters. In: *Practical Guide to Heavy Tails: Statistical Techniques and Applications* (R. J. Adler, R. E. Feldman and M. S. Taqqu, Eds.). pp. 311–335. Birkhäuser. Boston.

Kokoszka, P. S. and M. S. Taqqu (1994). Infinite variance stable ARMA processes. *Journal of Time Series Analysis* **15**, 203–220.

Koutrouvelis, I. A. (1980). A goodness-of-fit test of simple hypotheses based on the empirical characteristic function. *Biometrika* **67**, 238–240.

Lehmann, E.L. (1999). *Elements of Large-Sample Theory*. Springer Verlag. New York.

Nikias, C. L. and M. Shao (1995). *Signal Processing with α-Stable Distributions and Applications*. Wiley. New York.

Nowicka, J. (1997). Asymptotic behavior of the covariation and the codifference for ARMA models with stable innovations. *Stochastic Models* **13**, 673–686.

Nowicka, J. and A. Weron (1997). Measures of dependence for ARMA models with stable innovations. *Annales UMCS, Section A (Mathematica)* **51**, 133–44.

Proakis, J. G. and D. G. Manolakis (1996). *Digital Signal Processing*. Prentice Hall. London.

R Development Core Team (2004). *R: A language and environment for statistical computing*. R Foundation for Statistical Computing. Vienna, Austria. ISBN 3-900051-00-3.

Rachev, S. T. and S. Mittnik (2000). *Stable Paretian Models in Finance*. Wiley. New York.

Samorodnitsky, G. and M. S. Taqqu (1994). *Stable NonGaussian Processes: Stochastic Models with Infinite Variance*. Chapman and Hall. New York.

Yang, Q., A. Petropulu and J. C. Pesquet (2001). Estimating long-range dependence in impulsive traffic flows. In: *IEEE International Conference on Acoustics, Speech, and Signal Processing (ICASSP)*. pp. 3413 –3416.

Yu, J. (2004). Empirical characteristic function estimation and its applications. *Econometric Reviews* **23**, 93–123.

Table 1: The true values $I(\cdot)$ and the estimates $\hat{I}(\cdot)$

N	α	Method	s	$I(1)$	Avg. $\hat{I}(1)$	MAD_1	$I(2)$	Avg. $\hat{I}(2)$	MAD_2
100	2	Avg.	$\{0.01\}$		0.66180	0.04600		0.14647	0.10186
		Exp.	$\{0.01, 0.1, 1\}$	0.67722	0.63195	0.05732	0.17821	0.19387	0.09025
		ACF	-		0.65848	0.04644		0.14500	0.10115
	1.8	Avg.	$\{0.01, 0.2\}$	0.64700	0.64938	0.04415	0.19237	0.15860	0.09297
		Exp.	$\{0.01, 0.1, 1\}$		0.62509	0.04760		0.20409	0.08011
	1.5	Avg.	$\{0.01, 0.1, 0.2\}$	0.59903	0.60826	0.04839	0.21343	0.17158	0.07927
		Exp.	$\{0.01, 0.1, 1\}$		0.59728	0.04888		0.21062	0.06729
	1.3	Avg.	$\{0.01, 0.06, \ldots, 0.21\}$	0.56554	0.57235	0.05032	0.22665	0.19017	0.06974
		Exp.	$\{0.01, 0.1, 1\}$		0.57364	0.05548		0.22380	0.06105
	1	Avg.	$\{0.01, 0.02, \ldots, 0.2\}$	0.51350	0.50708	0.04745	0.24325	0.21576	0.06049
		Exp	$\{0.01, 0.02, \ldots, 0.2\}$		0.50717	0.04739		0.21577	0.06036
	0.8	Avg.	$\{0.01, 0.02, \ldots, 0.1\}$	0.47792	0.47274	0.04357	0.25014	0.22644	0.05599
		Exp.	$\{0.01, 0.02, \ldots, 0.1\}$		0.47276	0.04355		0.22643	0.05597
	0.5	Avg.	$\{0.01, 0.02, \ldots, 0.1\}$	0.42379	0.40713	0.04463	0.24809	0.22776	0.06033
		Exp.	$\{0.01, 0.02, \ldots, 0.1\}$		0.40715	0.04459		0.22776	0.06028
1000	2	Avg.	$\{0.01\}$		0.67577	0.01335		0.17495	0.03096
		Exp.	$\{0.01\}$	0.67722	0.67577	0.01335	0.17821	0.17495	0.03096
		ACF	-		0.67544	0.01336		0.17477	0.03094
	1.8	Avg.	$\{0.01, 0.06, \ldots, 0.21\}$	0.64700	0.65059	0.01645	0.19237	0.18859	0.02708
		Exp.	$\{0.01, 0.06, \ldots, 0.21\}$		0.65067	0.01649		0.18854	0.02708
	1.5	Avg.	$\{0.01, 0.02, \ldots, 0.2\}$	0.59903	0.60204	0.01815	0.21343	0.20879	0.02155
		Exp.	$\{0.01, 0.02, \ldots, 0.2\}$		0.60209	0.01822		0.20878	0.02154
	1.3	Avg.	$\{0.01, 0.02, \ldots, 0.2\}$	0.56554	0.56751	0.01640	0.22665	0.22312	0.01996
		Exp	$\{0.01, 0.02, \ldots, 0.2\}$		0.56754	0.01644		0.22310	0.01992
	1	Avg.	$\{0.01, 0.02, \ldots, 0.1\}$	0.51350	0.51396	0.01521	0.24325	0.24105	0.01577
		Exp.	$\{0.01, 0.02, \ldots, 0.1\}$		0.51396	0.01521		0.24105	0.01576
	0.8	Avg.	$\{0.01, 0.02, \ldots, 0.1\}$	0.47792	0.47742	0.01383	0.25014	0.24811	0.01511
		Exp.	$\{0.01, 0.02, \ldots, 0.1\}$		0.47742	0.01382		0.24811	0.01510
	0.5	Avg.	$\{0.01, 0.02, \ldots, 0.1\}$	0.42379	0.42239	0.01296	0.24809	0.24548	0.01852
		Exp.	$\{0.01, 0.02, \ldots, 0.1\}$		0.42239	0.01295		0.24548	0.01851

The true values $I(\cdot)$ and the estimates $\hat{I}(\cdot)$ from the experiment I, that is MA(2) process with $c_0 = 1$, $c_1 = 2$ and $c_2 = 1.111$ for $T = 1000$ replication, and for some sample size N. The ϵ_t is $S\alpha S$ process with some α and $\sigma = 1$. Here, $Avg.\hat{I}(i) = \frac{1}{T}\sum_{j=1}^{T} \operatorname{Re}\hat{I}(i)_j$, and $MAD_i = \frac{1}{T}\sum_{j=1}^{T} |\operatorname{Re}\hat{I}(i)_j - I(i)|$, $i = 1, 2$, where $\operatorname{Re}\hat{I}(i)_j$ denotes the estimates at lag i in run j . The weighting methods here denote by the simple average (method Avg.) and the negative exponential weighted average (method Exp.). Further explanation about the table is given in Section 3.2

Figure 1: Plots of Re $\hat{I}(1)$
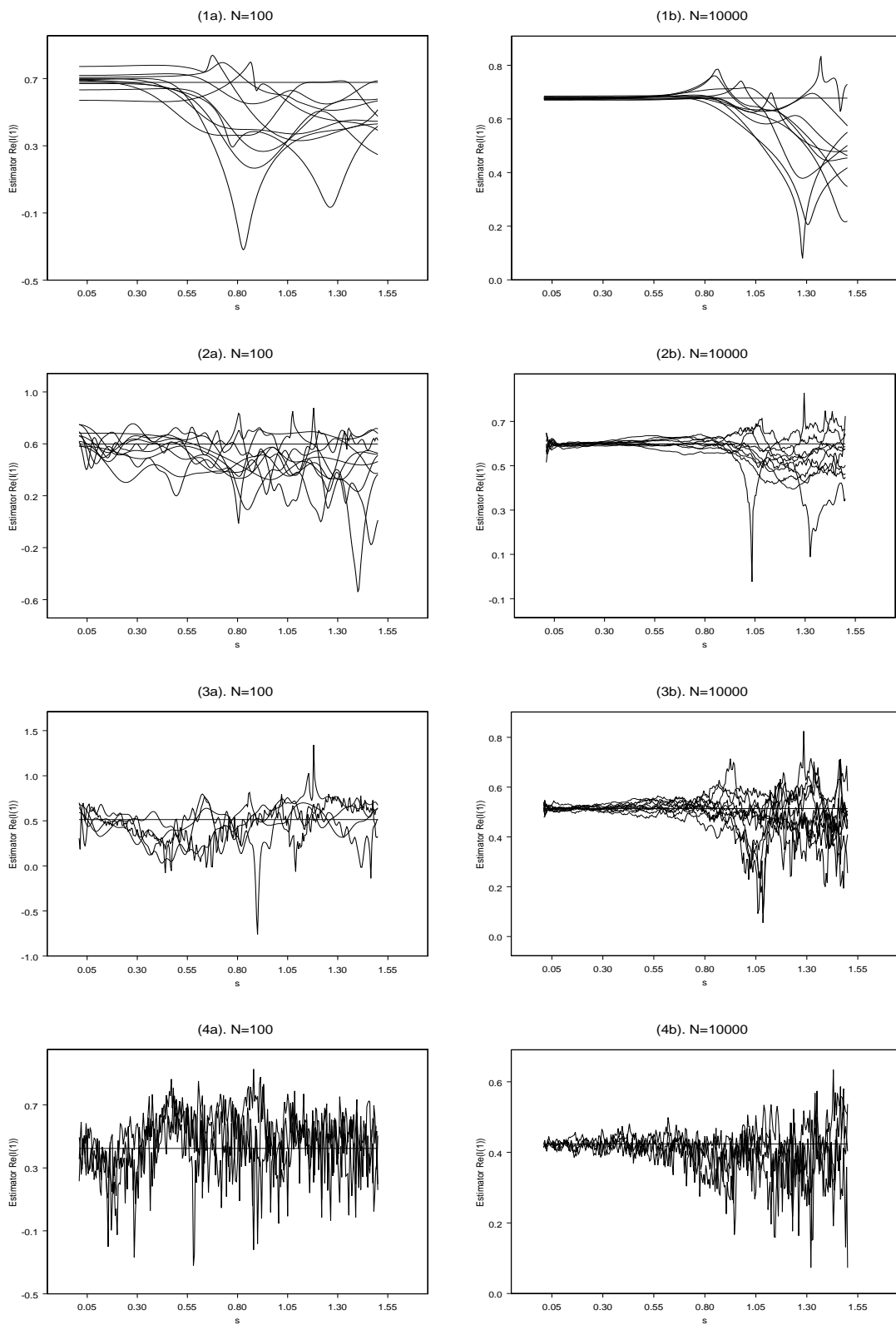
Plots of Re $\hat{I}(1)$ for several simulation runs where $(1a-1b).\alpha = 2$, $(2a-2b).\alpha = 1.5$, $(3a-3b).\alpha = 1$ and $(4a-4b).\alpha = 0.5$ and $\sigma = 1$, $s \in [0.01, 1.55]$. Data are generated from experiment I that is MA(2) process with $c_0 = 1$, $c_1 = 2$ and $c_2 = 1.111$. The straight lines denote the true values of $I(1)$.
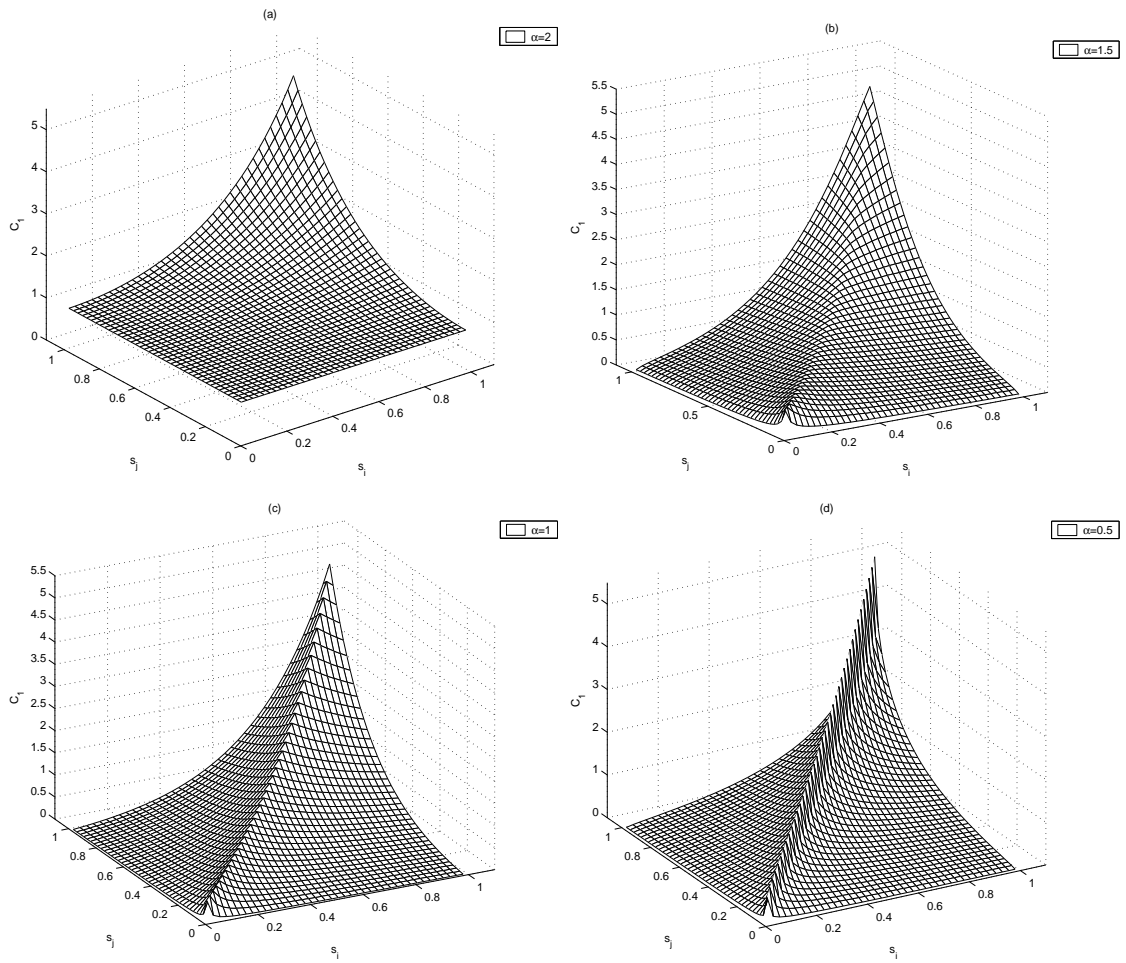
Figure 2: Plots of $W_1(i,j)$ (see eq. (16)), for $s_i, s_j \in [0.01, 1]$, and some α's.