

# Data-driven model selection for same-realization predictions in autoregressive processes

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#### **Abstract**

This paper is about the one-step ahead prediction of the future of observations drawn from an infinite-order autoregressive AR( $\infty$ ) process. It aims to design penalties (fully data driven) ensuring that the selected model verifies the efficiency property but in the non-asymptotic framework. We show that the excess risk of the selected estimator enjoys the best bias-variance trade-off over the considered collection. To achieve these results, we needed to overcome the dependence difficulties by following a classical approach which consists in restricting to a set where the empirical covariance matrix is equivalent to the theoretical one. We show that this event happens with probability larger than  $1-c_0/n^2$  with  $c_0>0$ . The proposed data-driven criteria are based on the minimization of the penalized criterion akin to the Mallows's  $C_p$ .

**Keywords** Model selection  $\cdot$  Oracle inequality  $\cdot$  Efficiency  $\cdot$  Autoregressive process  $\cdot$  Data driven

#### 1 Introduction

Consider observations  $(X_1, X_2, \dots, X_n)$  arising from a trajectory of the process

$$X_t = f^* \left( (X_{t-i})_{i \in \mathbb{N}^*} \right) + \sigma \, \xi_t \quad \text{for any } t \in \mathbb{Z}, \tag{1}$$

where  $(\xi_t)_{t\in\mathbb{Z}}$  is a sequence of zero-mean independent identically distributed random variables (i.i.d.r.v) satisfying  $\mathbb{E}(|\xi_0|^4)<\infty$  and  $f^*:\mathbb{R}^\mathbb{N}\to\mathbb{R}$  is a measurable function and  $\sigma>0$  an unknown constant.

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The problem is to estimate the function  $f^*$  using these observations. Process (1) is a particular case of the general class of affine causal process studied in Doukhan and Wintenberger (2008) and Bardet and Wintenberger (2009).

The study of this type of process more often requires the classical regularity condition on the function  $f^*$ , which are not restrictive at all and remain valid in various time series models. This condition can be stated as follows:

$$\sum_{k=1}^{\infty} \left( \sup_{x \in \mathbb{R}^{\infty}} \left| \frac{\partial}{\partial x_k} f^*(x) \right| \right) < 1, \tag{2}$$

provided that that  $f^*$  admits partial derivatives on  $\mathbb{R}^{\mathbb{N}}$ . Under (2) and if the noise  $\xi_0$  admits r-order moments, Doukhan and Wintenberger (2008) showed that there exists a stationary, mixing and ergodic solution to (1) admitting r-order moments.

Moreover, Bardet and Wintenberger (2009) studied the consistency and the asymptotic normality of the Quasi-maximum log-likelihood estimator (QMLE) of  $\theta^* = (\theta_i^*)_{i \in \mathbb{N}}$  in the case  $f^* = f_{\theta^*}$ .

In this paper, we will focus only on processes with a linear regression function  $(f_{\theta^*})$  with respect to the past and depending on some parameter  $\theta^* \in \mathbb{R}^{\mathbb{N}}$ ; that is

$$f^*(X_{t-1}, X_{t-2}, \dots) = f_{\theta^*}(X_{t-1}, X_{t-2}, \dots) = \sum_{i=1}^{\infty} \theta_i^* X_{t-i}.$$
 (3)

For such processes, condition (2) is rewritten as

$$\mathbf{A1}: \qquad \sum_{i=1}^{\infty} |\theta_i^*| < 1.$$

Even if this condition reduces the set of parameters a bit, the class of  $AR(\infty)$  processes checking the condition A1 is rich and of practical importance because it contains all invertible causal ARMA(p, q) processes and it is very useful for prediction given the past. Moreover, contrary to the autocovariance of ARMA(p, q) processes which decays exponentially fast,  $AR(\infty)$  are able to model more complex behavior such as slower decay of the covariance structure.

Henceforth, let observations  $(X_1, X_2, ..., X_n)$  be a trajectory of the solution  $X := (X_t)_{t \in \mathbb{Z}}$  of (1) verifying **A1**. The goal of this paper is to predict the next value  $X_{n+1}$ . In fact, if  $\theta^*$  were known, a simple prediction of  $X_{n+1}$  could be  $f_{\theta^*}(X_n, X_{n-1}, ...)$  setting  $X_t = 0$  for all  $t \le 0$ . However,  $\theta^*$  is generally unknown and it is impossible to provide a direct estimator since its coordinates are infinite. It is classical to identify a 'good' finite-dimensional model based on the data which can be done by sieve estimation where only a finite number of  $\{\theta_i^*\}_{i=1}^K$  is estimated and letting K grows as the sample size increases. A usual approach to this is model selection and the goal is to provide a model with the prediction error as small as the oracle's one.

This question has already been addressed in the literature. Shibata (1980) was the first to tackle this issue. He proved that Akaike criterion is *asymptotically efficient* in the sense that the selected model achieves a smaller one-step mean squared error of prediction when it is fitted to predict an independent realization of the same process.



Following Shibata's asymptotically setting, Ing and Wei (2003) and Ing and Wei (2005) extended this result for same realization predictions. Indeed, they argued that the Shibata's idea to fit the model to another independent realization is unrealistic since in practice we only have one data at hand. The common feature of these works is their asymptotic framework.

Meanwhile, there were several authors which study this question in non-asymptotic regime. Goldenshluger and Zeevi (2001) in the nonparametric framework, studied how well a Gaussian process admitting an  $AR(\infty)$  representation can be approximated by a finite-order AR model.

In Baraud et al. (2001a) and (b), they analyzed similar question, but a little bit different as observations arise from an auto-regressive model of order k. They proved an oracle inequality under several conditions, for instance the compactly supported base of the regression function. Moreover, they assume that the process is  $\beta$ -mixing, which is usually admitted, but quite hard to verify in practice. For linear processes, the  $\tau$ -mixing is more suitable since its coefficients can be easily computed (see Comte et al., 2008) and be bounded by a function of the model parameter  $\theta^*$  (see Doukhan and Wintenberger, 2008). In this work, we do not assume any mixing property of the process since the condition **A1** implies the  $\tau$ -mixing property (see Doukhan and Wintenberger, 2008) and we will see that the decreasing rate of  $\tau$ -mixing coefficients is bounded by the decreasing rate of the coefficients  $\theta^* = (\theta^*_i)_{i \in \mathbb{N}}$ .

Based on the above and following a model selection approach, our purpose in this work is to design adaptive penalties in such a way that the selected model mimics the oracle when observations arise from  $AR(\infty)$  under mild conditions, including the existence of the all order moments of the noise, the decreasing rate of the coefficients of  $(\theta_i^*)_{i \in \mathbb{N}}$  so that thanks to a result by Doukhan and Wintenberger (2008), the generating process has nice properties such as stationarity,  $\tau$ -mixing.

The main contribution of this paper is to prove that the excess risk of the selected estimator enjoys the best bias-variance trade-off over the considered collection.

The paper is organized as follows. The model selection approach along with preliminary results is described in Sect. 2. The main results are presented in Sect. 3. Finally, Sect. 4 contains the proofs.

# 2 Model selection approach and preliminary results

Before entering properly into the description of our approach, let us introduce some notations.

#### 2.1 Notations

We will use the following norms:

- < .,. > is the usual scalar product and if  $x, y \in \mathbb{R}^n$ , < .,. ><sub>n</sub> =  $\frac{1}{n} \sum_{i=1}^n x_i y_i$ ;
- The usual Euclidean norm on  $\mathbb{R}^{\nu}$ , with  $\nu \geq 1$ , is denoted by  $\|.\|^{\mu}$  and its normalized version by  $\|.\|_n$ ;



 $||A||_{op}$  is the operator norm of A as the square root of the largest eigenvalue of  $A^{\mathsf{T}}A$ . If A is symmetric, then  $||A||_{\mathrm{op}}$  is the largest (in absolute value) eigenvalue of

If X is a  $\mathbb{R}^{\nu}$ -random variable and  $r \ge 1$ , we set  $||X||_r = (\mathbb{E}[||X||^r])^{1/r} \in [0, \infty]$ , where  $\mathbb{E}[Y]$  denotes the expected value of the random variable Y.

#### 2.2 Model selection approach

Let  $S_m$  (shortly m) a model for  $f^*$  be the set of linear function f from  $\mathbb{R}^{D_m}$  to  $\mathbb{R}$  such

$$f(x_1, x_2, \dots, x_{D_m}) = f_{\theta}(x_1, x_2, \dots, x_{D_m}) = \sum_{i=1}^{D_m} \theta_i x_i,$$
 (4)

with  $\theta = (\theta_1, \dots, \theta_{D_m}) \in \Theta_m$  and  $\Theta_m$  a compact set of  $\mathbb{R}^{D_m}$  $\sup_{\theta \in \Theta_m} \sum_{i=1}^{D_m} |\theta_i| < 1.$ 

 $S_m$  can be viewed as an AR( $D_m$ ) model.

Given a predictor  $f_{\theta} \in S_m$ , its quality is measured by the quadratic loss

$$R(\theta) = \mathbb{E}\left[ (X_{n+1} - f_{\theta}^{n+1})^2 \right]$$

where  $f_{\theta}^{n} = f_{\theta}(X_{n-1}, \dots, X_{n-D_{m}})$ . According to Bardet and Wintenberger (2009), the Bayes predictor which minimizes  $R(\theta)$  over the set of all predictors is the inaccessible function  $f_{\theta^*}$ . Let then introduce the excess loss of the predictor  $f_{\theta}$  (with respect to  $f_{\theta^*}$ )

$$\mathcal{E}(\theta,\theta^*) := R(\theta) - R(\theta^*) = \mathbb{E}\left[ (f_{\theta^*}^{n+1} - f_{\theta}^{n+1})^2 \right] \ge 0.$$

Given a model m, we define its best predictor  $f_{\theta_m^*}$  by

$$\theta_m^* = \underset{\theta \in \Theta_m}{\operatorname{argmin}} R(\theta).$$

Its empirical version minimizing the least-squares contrast is

$$\widehat{\theta}_m = \underset{\theta \in \Theta_m}{\operatorname{argmin}} \ \gamma_n(\theta) \quad \text{where} \quad \gamma_n(\theta) = \frac{1}{n} \sum_{t=1}^n (X_t - f_\theta^t)^2. \tag{5}$$

In this work, we will consider that the excess loss is measured on the design points, that is to say

$$\mathscr{E}(\widehat{\theta}, \theta^*) = \mathbb{E}\left[\left\|F_{\widehat{\theta}} - F_{\theta^*}\right\|_n^2\right] \tag{6}$$

where  $F_{\theta} := (f_{\theta}^1, \dots, f_{\theta}^n)^{\mathsf{T}}$  and  $\|x\|_n^2 = \frac{1}{n} \sum_{t=1}^n x_t^2$ . Given that all the models which can be considered must have finite dimensions for fixed n, making all  $S_m$  wrong models, it is classical to let the dimension of



competitive models grow with the number of observations. This will help reduce the excess loss and provide a better approximation of  $f_{\theta^*}$ .

Let  $\mathcal{M}_n$  a countable collection of hierarchical model  $S_m$  and  $K_n$  is the dimension of the largest model in  $\mathcal{M}_n$  satisfying  $|\mathcal{M}_n| \leq K_n < n$ . We follow the classical approach of model selection which consists in minimizing the penalized LSE.

Let pen:  $\mathcal{M}_n \to \mathbb{R}^+$  be a penalty function, possibly data-dependent, and define

$$\widehat{m} = \underset{m \in \mathcal{M}_n}{\operatorname{argmin}} \{C(m)\} \quad \text{with} \quad C(m) := \gamma_n(\widehat{\theta}_m) + \operatorname{pen}(S_m).$$
 (7)

The function pen can be a linear function of the model dimension (see for instance Birgé and Massart, (2001, 2007) among others) or a non linear one (see Lebarbier and Mary-Huard, 2004).

The best possible choice over  $\mathcal{M}_n$  is  $m^*$  the so-called *oracle* defined as

$$m^* \in \arg\inf_{m \in \mathcal{M}_n} \ell(\widehat{\theta}_m, \theta^*).$$
 (8)

The oracle  $m^*$  is unachievable since it depends on  $\theta^*$  and the distribution  $P_{(X_1,...,X_n)}$  that are unknowns. However, we hope to select a model  $\widehat{m}$  so that  $\ell(\widehat{\theta}_{\widehat{m}},\theta^*)$  is closest to  $\ell(\widehat{\theta}_{m^*},\theta^*)$ .

The goal of this paper is to propose a data-driven penalty in order to obtain an oracle inequality

$$\ell(\widehat{\theta}_{\widehat{m}}, \theta^*) \le C_1 \inf_{m \in \mathcal{M}_n} \left\{ \ell(\widehat{\theta}_m, \theta^*) \right\} + \frac{C_2}{n} \tag{9}$$

with the leading constant  $C_1$  close to one and  $C_2 > 0$ . This goal could rather be to show that the excess risk of the selected estimator  $\hat{\theta}_{\widehat{m}}$  realizes the best bias-variance trade-off, which would make our penalty an ideal choice in terms of excess risk.

$$\mathscr{E}(\widehat{\theta}_{\widehat{m}}, \theta^*) \le C_1' \inf_{m \in \mathcal{M}_n} \left\{ \mathscr{E}(\theta_m^*, \theta^*) + \operatorname{pen}(S_m) \right\} + \frac{C_2'}{n}$$
(10)

with the leading constant  $C_1' = 1 + \delta$  with  $\delta > 0$  (and close to 0) and  $C_2' > 0$ .

That is to say that the selected model  $\hat{m}$  will be large enough to reduce its bias, but not too large to avoid high variance.

#### 2.3 Preliminary results and assumptions

As we are in dependence setting, we are going to leverage the  $\tau$ -mixing property of  $(X_t)_{t \in \mathbb{Z}}$  in order to obtain some exponential inequalities. The  $\tau$ -mixing coefficients are a measure of the dependence of the process and has been introduced by Dedecker and Prieur (2005). This will help us build 'independents' random vectors and apply classical exponential inequalities. Let then introduce some notations.

Let  $(\Omega, \mathcal{C}, \mathbb{P})$  be a probability space,  $\mathcal{M}$  a  $\sigma$ -subalgebra of  $\mathcal{C}$  and Z a random variable with values in a Banach space  $(E, \|.\|_E)$ . Assume that  $\mathbb{E}|Z| < \infty$  and define



$$\tau^{(p)}(\mathcal{M}, Z) = \left\| \sup_{f \in \Lambda(E)} \left\{ \left| \int f(x) \mathbb{P}_{Z|\mathcal{M}}(dx) - \int f(x) \mathbb{P}_{Z}(dx) \right| \right\} \right\|_{p}$$

where  $\Lambda(E)$  is the set of 1-Lipschitz function, i.e., the functions f from  $(E, ||.||_E)$  to  $\mathbb{R}$  such that  $|f(x) - f(y)| \le ||x - y||_E$ .

Using definition of  $\tau$ , we will measure the dependence of the strictly stationary sequence  $(Z_t)_{t\in\mathbb{Z}}$  thanks to the coefficients defined as follows. For any  $s\geq 0$ , let introduce the norm  $\|x-y\|_{\mathbb{R}^k}=(|x_1-y_1|+\cdots+|x_k-y_k|)$  and setting  $\mathcal{M}_i=\sigma(Z_t,t\leq i)$  and if  $\mathbb{E}(|Z_1|)<\infty$ , let

$$\tau_{Z,\infty}^{(p)}(s) = \sup_{l>0} \ \Big\{ \max_{1\leq k\leq l} \frac{1}{k} \ \sup\Big\{ \tau^{(p)} \Big(\mathcal{M}_i, (Z_{i_1},\dots,Z_{i_k}) \Big) \ , i+s \leq i_1 < \dots < i_k \Big\} \Big\}.$$

Finally, the time series  $(Z_t)_{t \in \mathbb{Z}}$  is  $\tau_{Z,\infty}^{(p)}$ -weakly dependent when its coefficients  $\tau_{Z,\infty}^{(p)}$  tend to 0 as s tends to infinity.

Proposition 3 that is a consequence of Theorem 3.1 in Doukhan and Wintenberger (2008) gives a link between the  $\tau$ -mixing coefficients of the process  $(X_t)_{t \in \mathbb{Z}}$  and the coefficients  $\theta_i^*$  of model (3).

As we are going to need independence for block of random variables, let denote for t = 1, ..., n the random vector  $\mathbf{X}_t := (X_{t-1}, ..., X_{t-K_n})^{\mathsf{T}}$  One can see that the process  $(\mathbf{X}_t)_{t \in \mathbb{Z}}$  is also mixing with  $\tau_{\mathbf{X}_{\infty}}^{(1)}$  upper bounded by  $K_n \tau_{X_{\infty}}^{(1)}$  (see Lemma 1).

Now, we construct random variables approximating  $X_t$ 's enjoying the independence by block property. Let  $s_n$ ,  $q_n$  two integers such that  $n = 2 s_n q_n$ . We are going to build  $2 s_n$  blocks of length  $q_n$  so that the even index blocks are independent and so the odd index blocks.

For  $k = 0, \dots, s_n - 1$  let denote by

$$A_k = (\mathbf{X}_{2kq_n+1}, \dots, \mathbf{X}_{(2k+1)q_n})$$
 and  $B_k = (\mathbf{X}_{(2k+1)q_n+1}, \dots, \mathbf{X}_{(2k+2)q_n})$ .

Proposition 4 recalls a result from Lerasle (2011) that is a consequence of the coupling in Dedecker and Prieur (2005). It allows to have the block independence property.

To prove the oracle inequality, we will assume some constraints on the observations. **A2**  $X_t$  is sub-Gaussian with variance proxy  $\sigma_0^2 > 0$  i.e.,

$$\mathbb{E}[e^{\lambda X_t}] \le e^{\lambda^2 \sigma_0^2/2} \quad \text{for any } \lambda > 0.$$

Condition A2 implies that the vector  $Z_t^m = (X_{t-1}, \dots, X_{t-D_m})^{\mathsf{T}}$  which will be prominent in the proofs, is sub-Gaussian with variance proxy  $D_m \, \sigma_0^2$ . Indeed for any  $v \in \mathbb{R}^{D_m}$  such that  $\|v\| = 1$ ,



$$\begin{split} \mathbb{E}\Big[\exp\left(\lambda\,v^{\mathsf{T}}Z_{t}^{m}\right)\Big] &= \mathbb{E}\bigg[\prod_{i=1}^{D_{m}}\exp\left(\lambda\,v_{i}X_{t-i}\right)\bigg] \\ \≤\prod_{i=1}^{D_{m}} \left\|\exp\left(\lambda\,v_{i}X_{t-i}\right)\right\|_{D_{m}} \\ &\leq \prod_{i=1}^{D_{m}}\exp\left(\lambda^{2}\,D_{m}\,\sigma_{0}^{2}\,v_{i}^{2}/2\right) \\ &\leq e^{\frac{\lambda^{2}}{2}\,D_{m}\,\sigma_{0}^{2}}, \end{split}$$

where the Inequality follows from Hölder's Inequality.

The following assumption provides a sufficient condition to ensure the invertibility of both  $\widehat{\Sigma}_m := \mathbf{M}_m^{\mathsf{T}} \mathbf{M}_m$  and  $\Sigma_m := \mathbb{E} \left[ \widehat{\Sigma}_m \right]$  where  $\mathbf{M}_m = \left[ X_{i-1}, \dots, X_{i-D_m} \right]_{i=1}^n$ .

**A3**: For any  $f_{\theta} \in S_m$ ,  $\langle \alpha, \partial_{\theta} f_{\theta} \rangle = 0$  a.s.  $\Longrightarrow \alpha = 0$ 

This condition means that the columns of the matrix  $\mathbf{M}_m$  are linearly independents.

We will also need to bound eigenvalues of the matrices  $\Sigma_m$  for any  $m \in \mathcal{M}_n$ . To do that, we will leverage the relation between the spectral density of the process and these eigenvalues. Let us denote by r, the covariance function  $r(h) := \mathbb{E}[X_t X_{t+h}]$  for any integer h. Let also introduce the function  $g : [-\pi, \pi[ \longrightarrow \mathbb{C}$  such that for any  $\lambda$ ,

$$g(\lambda) = \frac{1}{2\pi} \sum_{h \in \mathbb{Z}} r(h) e^{-ih\lambda},$$

which exists under **A1** with  $|\theta_t^*| = O(t^{-\gamma})$  where  $\gamma \ge 1$ . Therefore, r is the inverse transform of g and  $r(h) = \int_{-\pi}^{\pi} e^{ih\lambda} g(\lambda) d\lambda$  for any  $h \in \mathbb{Z}$ . We will assume that

**A4**: There exists a constant a > 0 such that  $\inf_{-\pi \le \lambda < \pi} g(\lambda) \ge a$ .

This is a very weak assumption, and we are going to give the value of a for AR(p) process with  $p \in \mathbb{N}$  and  $p \ge 1$ . Let denote  $\theta^*(z) = 1 - \sum_{j=1}^p \theta_j^* z^j$ , it is well known for such process that

$$g(\lambda) = \frac{\sigma^2}{2 \pi |\theta^*(e^{-i\lambda})|^2}.$$

For instance for p equal to one, and  $X_t = \theta_1^* X_{t-1} + \sigma \xi_t$  with  $|\theta_1^*| < 1$ , it follows

$$g(\lambda) = \frac{\sigma^2}{2\pi \left| 1 - \theta_1^* e^{-i\lambda} \right|^2} = \frac{\sigma^2}{2\pi \left( 1 - 2 \theta_1^* \cos(\lambda) + (\theta_1^*)^2 \right)},$$

and then it is simple to see that



$$a := \frac{\sigma^2}{2\pi (1 + |\theta_1^*|)^2} \le g(\lambda) \le \frac{\sigma^2}{2\pi (1 - |\theta_1^*|)^2}.$$

For  $p \ge 1$  and  $X_t = \sum_{j=1}^p \theta_j^* X_{t-j} + \sigma \xi_t$  satisfying  $\sum_{j=1}^p \theta_j^* < 1$  and  $\theta_j^* \ge 0$ , we have

$$g(\lambda) = \frac{\sigma^2}{2\pi \left| 1 - \sum_{j=1}^p \theta_j^* e^{-ij\lambda} \right|^2}$$

$$= \sigma^2 (2\pi)^{-1} \left( 1 + \sum_{j=1}^p (\theta_j^*)^2 - 2 \sum_{j=1}^p \theta_j^* \cos(j\lambda) + 2 \sum_{k=1}^{p-1} \theta_k^* \left\{ \sum_{j=k+1}^p \theta_j^* \cos\left((j-k)\lambda\right) \right\} \right)^{-1}.$$

Thus, using  $|\cos(x)| \le 1$  for any real x, it follows for every  $\lambda$  that

$$\sigma^{2} (2\pi)^{-1} \left( 1 + \sum_{j=1}^{p} (\theta_{j}^{*})^{2} + 2 \sum_{j=1}^{p} \theta_{j}^{*} + 2 \sum_{k=1}^{p-1} \theta_{k}^{*} \left\{ \sum_{j=k+1}^{p} \theta_{j}^{*} \right\} \right)^{-1} \leq g(\lambda).$$

For such AR(p) process, one can take the constant a in A4 to be equal to

$$a = \sigma^2 (2\pi)^{-1} \left( 1 + \sum_{j=1}^p (\theta_j^*)^2 + 2 \sum_{j=1}^p \theta_j^* + 2 \sum_{k=1}^{p-1} \theta_k^* \left\{ \sum_{j=k+1}^p \theta_j^* \right\} \right)^{-1}.$$

We can now state an important intermediate result which provides uniform lower and upper bound on the spectral norm of the matrices  $\Sigma_m$ .

**Proposition 1** Under A1 with  $|\theta_t^*| = O(t^{-\gamma})$  where  $\gamma \ge 2$ , we have for any  $m \in \mathcal{M}_n$ 

$$\|\Sigma_m\|_{\text{op}} \le \pi^{-1} \sum_{i=0}^{\infty} |\mathbb{E}[X_0 X_i]| < \infty.$$
 (11)

Moreover and under A3-A4, it holds

$$\left\| \Sigma_m^{-1} \right\|_{\text{on}} \le 1/a. \tag{12}$$

Let us introduce extra important notations. Let denote by  $\mu$  the law of the vector  $\mathbf{X}_t$  and

$$\Omega_n = \left\{ \omega : \left| \frac{\|F_\theta\|_n^2}{\|F_\theta\|_\mu^2} - 1 \right| \le \frac{1}{2}, \quad \forall F_\theta \in \bigcup_{m,m' \in \mathcal{M}_n} (S_m + S_{m'}) \right\}$$

where  $\|F_{\theta}\|_{\mu}^2 := \frac{1}{n} \mathbb{E} \left[ \sum_{t=1}^n (f_{\theta}^t)^2 \right] = \int (f_{\theta}^1)^2 d\mu$ . It is common to consider the set  $\Omega_n$  which makes a link between the empirical norm  $\|.\|_n$  and the to  $\mathbb{L}_2$  norm (see for instance Baraud et al., 2001b; Hsu et al., 2011; Van de Geer, 2002, Comte & Genon-Catalot, 2020 among others). We will see that in our framework,  $\Omega_n$  holds with large probability.

In all of this work, we assume that  $q_n$  was chosen to verify



$$\mathbf{A5}: \qquad \left(\frac{\log q_n}{q_n}\right)^{\gamma-1} \le \frac{A}{n},\tag{13}$$

for some constant A and  $\gamma > 1$ . Also we choose the integer  $s_n$  such that

**A6**: 
$$\frac{s_n}{2} \min \left\{ \left( \frac{1}{2^7 \sigma_0^2 K_n} \right)^2, \frac{1}{2^8 \sigma_0^2 K_n} \right\} \ge 3 \log n.$$
 (14)

This means that  $s_n$  is of the form  $s_n = C \log n$  where  $c \ge 6 \max \left\{ (2^7 \sigma_0^2 K_n)^2, 2^8 \sigma_0^2 K_n \right\}$ .

**Proposition 2** *Under assumptions* **A1**, **A6** *and if*  $|\theta_t^*| = O(t^{-\gamma})$  *with*  $\gamma \ge 8$ , *there exists a constant C such that* 

$$\mathbb{P}(\Omega_n^c) \le \frac{C}{n^3}.\tag{15}$$

#### 3 Bias-variance result

We are now able to state the main result of the paper.

**Theorem 1** Let consider observations  $(X_1, \ldots, X_n)$  arising from a solution of process (1) satisfying **A1** with  $|\theta_t^*| = O(t^{-\gamma})$  where  $\gamma \geq 8$  and also verifying **A2** and **A4**. Let  $\mathcal{M}_n$  be some countable family of AR models satisfying **A3** and **A5-A6**. For any constant x > 2, let a penalty function pen:  $\mathcal{M}_n \to \mathbb{R}^+$  such that

$$pen(S_m) \ge 8x^3\sigma^2 \frac{D_m}{n}.$$
 (16)

Then, the LSE  $\hat{\theta}_{\hat{m}}$  with  $\hat{m}$  given in (7), satisfies

$$\mathbb{E}\Big[ \|F_{\hat{\theta}_{\hat{m}}} - F_{\theta^*}\|_n^2 \Pi_{\Omega_n} \Big] \le C_1(x) \inf_{m \in \mathcal{M}_n} \left\{ \mathbb{E}\Big[ \|F_{\theta^*} - F_{\theta_m^*}\|_n^2 \Big] + 2 \operatorname{pen}(S_m) \right\} + \frac{x(x+2)}{x-2} \frac{C_2}{n}$$
(17)

where  $C_1(x) = \left(\frac{x+2}{x-2}\right)^2 > 1$  and  $C_2 > 0$ .

As we can see, this result is almost similar to that of Baraud et al. (2001b) obtained in nonparametric framework. However, their result is only valid if we want to estimate the function  $F_{\theta^*}$  on some compact set. This restriction is lifted in our parametric framework.



## 4 Proofs

#### 4.1 Proof of Theorem 1

**Proof** We follow the scheme of the proof of Baraud et al. (2001b). Let fix  $m \in \mathcal{M}_n$ . From definition (7), we have

$$\gamma_n(\widehat{\theta}_{\widehat{m}}) + \operatorname{pen}(S_{\widehat{m}}) \le \gamma_n(\widehat{\theta}_m) + \operatorname{pen}(S_m) \le \gamma_n(\theta_m^*) + \operatorname{pen}(S_m).$$
 (18)

Since,

$$\gamma_n \left( \widehat{\theta}_m \right) = \frac{1}{n} \sum_{t=1}^n (X_t - f_{\widehat{\theta}_m}^t)^2 = \frac{1}{n} \sum_{t=1}^n (f_{\theta^*}^t - f_{\widehat{\theta}_m}^t)^2 + \frac{\sigma^2}{n} \sum_{t=1}^n \xi_t^2 - \frac{2\sigma}{n} \sum_{t=1}^n \xi_t (f_{\widehat{\theta}_m}^t - f_{\theta^*}^t),$$

(18) yields to

$$\|F_{\hat{\theta}_{\hat{m}}} - F_{\theta^*}\|_n^2 \le \|F_{\theta^*_m} - F_{\theta^*}\|_n^2 + \frac{2\sigma}{n} \sum_{t=1}^n \xi_t (f_{\hat{\theta}_{\hat{m}}}^t - f_{\theta^*_m}^t) + \operatorname{pen}(S_m) - \operatorname{pen}(S_{\hat{m}}).$$
(19)

The difficult part of this proof is to handle the inner product  $\frac{2}{n} \sum_{t=1}^{n} \sigma \, \xi_{t} (f_{\widehat{\theta}_{\widehat{m}}}^{t} - f_{\theta_{m}^{*}}^{t})$ , which can be rewritten as

$$\begin{split} \frac{2}{n} \sum_{t=1}^{n} \sigma \, \xi_{t} (f_{\widehat{\theta}_{\widehat{m}}}^{t} - f_{\theta_{m}^{*}}^{t}) &= \frac{2}{n} \| F_{\widehat{\theta}_{\widehat{m}}} - F_{\theta_{m}^{*}} \|_{\mu} \sum_{t=1}^{n} \sigma \, \xi_{t} \frac{(f_{\widehat{\theta}_{\widehat{m}}}^{t} - f_{\theta_{m}^{*}}^{t})}{\| F_{\widehat{\theta}_{\widehat{m}}} - F_{\theta_{m}^{*}} \|_{\mu}} \\ &\leq \frac{2}{n} \| F_{\widehat{\theta}_{\widehat{m}}} - F_{\theta_{m}^{*}} \|_{\mu} \sup_{g \in B(\widehat{m}, \mu)} \sum_{t=1}^{n} \sigma \, \xi_{t} \, g_{\theta}^{t} \\ &\leq x^{-1} \| F_{\widehat{\theta}_{\widehat{m}}} - F_{\theta_{m}^{*}} \|_{\mu}^{2} + n^{-2} \, x \left( \sup_{g \in B(\widehat{m}, \mu)} \sum_{t=1}^{n} \sigma \, \xi_{t} \, g_{\theta}^{t} \right)^{2} \end{split}$$

since  $2 a b \le x^{-1} a^2 + x b^2$  for any x > 0 and where

$$B(m', \mu) = \left\{ F_{\theta} \in S_m + S_{m'} : \|F_{\theta}\|_{\mu}^2 \le 1 \right\}.$$

Moreover, on the set  $\Omega_n$ , it holds

$$\begin{split} \left\| F_{\widehat{\theta}_{\widehat{m}}} - F_{\theta_{m}^{*}} \right\|_{\mu}^{2} &\leq 2 \left\| F_{\widehat{\theta}_{\widehat{m}}} - F_{\theta_{m}^{*}} \right\|_{n}^{2} \\ &\leq 2 \left( \left\| F_{\widehat{\theta}_{\widehat{m}}} - F_{\theta^{*}} \right\|_{n} + \left\| F_{\theta^{*}} - F_{\theta_{m}^{*}} \right\|_{n} \right)^{2} \\ &\leq 2 \left( 1 + y \right) \left\| F_{\widehat{\theta}_{\widehat{n}}} - F_{\theta^{*}} \right\|_{n}^{2} + 2 \left( 1 + y^{-1} \right) \left\| F_{\theta^{*}} - F_{\theta_{m}^{*}} \right\|_{n}^{2} \end{split}$$

for some y > 0. As a result,



$$\begin{split} & \frac{2}{n} \sum_{t=1}^{n} \sigma \, \xi_{t} (f_{\widehat{\theta}_{\widehat{m}}}^{t} - f_{\theta_{m}^{*}}^{t}) \leq 2 \, \frac{(1+y)}{x} \, \left\| F_{\widehat{\theta}_{\widehat{m}}} - F_{\theta^{*}} \right\|_{n}^{2} + 2 \, \frac{(1+y^{-1})}{x} \, \left\| F_{\theta^{*}} - f_{\theta_{m}^{*}} \right\|_{n}^{2} \\ & + \frac{x}{n^{2}} \, \left( \sup_{g \in B(\widehat{m}, \mu)} \sum_{t=1}^{n} \sigma \, \xi_{t} \, g_{\theta}^{t} \right)^{2}. \end{split}$$

Therefore, from (19), it holds on  $\Omega_n$ 

$$\begin{split} & \left(1 - 2\frac{(1+y)}{x}\right) \left\|F_{\hat{\theta}_{\hat{m}}} - F_{\theta^*}\right\|_n^2 \leq \left(1 + 2\frac{(1+y^{-1})}{x}\right) \left\|F_{\theta^*} - F_{\theta^*_m}\right\|_n^2 \\ & + \operatorname{pen}(S_m) - \operatorname{pen}(S_{\hat{m}}) + \frac{x\,\sigma^2}{n^2} \left(\sup_{g \in B(\hat{m},\mu)} \sum_{t=1}^n \xi_t \, g^t_{\theta}\right)^2 \\ & \leq \left(1 + 2\frac{(1+y^{-1})}{x}\right) \left\|F_{\theta^*} - F_{\theta^*_m}\right\|_n^2 + \operatorname{pen}(S_m) - \operatorname{pen}(S_{\hat{m}}) \\ & + 8\,x^3\sigma^2\, \frac{D_m + D_{\hat{m}}}{n} + \frac{x\,\sigma^2}{n^2} \left[\left(\sup_{g \in B(\hat{m},\mu)} \sum_{t=1}^n \, \xi_t \, g^t_{\theta}\right)^2 - 8n\,x^2 D(S_{\hat{m}})\right]_+ \end{split}$$

where  $D(S_{\widehat{m}}) = dim(S_m + S_{\widehat{m}}) \le D_m + D_{\widehat{m}}$ . Hence, using the condition on penalty (16),

$$\left(1 - 2\frac{(1+y)}{x}\right) \left\|F_{\widehat{\theta}_{\widehat{m}}} - F_{\theta^*}\right\|_n^2 \le \left(1 + 2\frac{(1+y^{-1})}{x}\right) \left\|F_{\theta^*} - F_{\theta_m^*}\right\|_n^2 + 2\operatorname{pen}(S_m) + x\,\sigma^2\,V_{\widehat{m}} \tag{20}$$

with

$$V_{m'} = \left[ \left( \sup_{g_{\theta} \in B(m', \mu)} v_n(g_{\theta}) \right)^2 - 8 \frac{x^2}{n} D(S_{m'}) \right]_+,$$

where  $v_n(g_\theta) := n^{-1} \sum_{t=1}^n \xi_t g_\theta^t$ . The proof will be established after controlling the expectation of  $V_{\widehat{m}}$  which involves the supremum of an empirical process.

Now we leverage the mixing property in order to apply Talagrand's Inequality (Theorem 2) to tackle  $\mathbb{E}[V_{\widehat{m}}]$ .

We have

$$V_{\widehat{m}} = \left[ \sup_{g_{\theta} \in B(\widehat{m}, \mu)} \left( v_n(g_{\theta}) \right)^2 - 8 \frac{x^2}{n} D(S_{\widehat{m}}) \right]_{+} \le 2 \sup_{g_{\theta} \in B(\widehat{m}, \mu)} \left( v_n(g_{\theta}) - v_n^*(g_{\theta}) \right)^2 + V_{\widehat{m}}^*$$
(21)

where

$$V_{\widehat{m}}^* = \left[ 2 \sup_{g_{\theta} \in B(\widehat{m}, \mu)} \left( v_n^*(g_{\theta}) \right)^2 - 8 \frac{x^2}{n} D(S_{\widehat{m}}) \right]_{+} \quad \text{and} \quad v_n^*(g_{\theta}) = \frac{1}{n} \sum_{t=1}^n \xi_t g_{\theta}(\mathbf{X}_t^*).$$



1/ **Control of**  $\mathbb{E}\Big[\sup_{g_{\theta}\in B(\widehat{m},\mu)}\Big(v_n(g_{\theta})-v_n^*(g_{\theta})\Big)^2\Big]$ . Let  $m'\in\mathcal{M}_n$  and  $g_{\theta}\in B(m',\mu)$ . Since the parameter set are compacts and  $\theta\mapsto g_{\theta}$  is continuous, there exists  $\theta_0\in\Theta_m\cup\Theta_{m'}$  such that

$$\sup_{g_{\theta} \in B(m',\mu)} \biggl( \nu_n(g_{\theta}) - \nu_n^*(g_{\theta}) \biggr)^2 = \frac{1}{n^2} \biggl( \ \sum_{t=1}^n \xi_t \left( g_{\theta_0}(\mathbf{X}_t) - g_{\theta_0}(\mathbf{X}_t^*) \right) \biggr)^2.$$

As  $\xi_t$  and  $\mathcal{F}_t$  are independents, it follows that

$$\mathbb{E}\left[\sup_{g_{\theta} \in B(m',\mu)} \left(v_n(g_{\theta}) - v_n^*(g_{\theta})\right)^2\right] = \frac{1}{n^2} \sum_{t=1}^n \mathbb{E}\left[\xi_t^2 \left(g_{\theta_0}(\mathbf{X}_t) - g_{\theta_0}(\mathbf{X}_t^*)\right)^2\right]$$
$$= \frac{1}{n} \mathbb{E}\left[\left(g_{\theta_0}(\mathbf{X}_0) - g_{\theta_0}(\mathbf{X}_0^*)\right)^2\right]$$

since  $\mathbb{E}[\xi_0^2] = 1$ . In addition,

$$\mathbb{E}\left[\left(g_{\theta_{0}}(\mathbf{X}_{0}) - g_{\theta_{0}}(\mathbf{X}_{0}^{*})\right)^{2}\right] = \sum_{i=1}^{D(S_{m'})} \sum_{j=1}^{D(S_{m'})} \theta_{0,i} \,\theta_{0,j} \,\mathbb{E}\left[(X_{-i} - X_{-i}^{*})(X_{-j} - X_{-j}^{*})\right] \\ \leq \left(\sum_{i=1}^{D(S_{m'})} \theta_{0,i} \,\|X_{-i} - X_{-i}^{*}\|_{2}\right)^{2}$$

using Cauchy-Schwarz Inequality. It then follows as  $\sum_{i=1}^{D(S_{m'})} |\theta_{0,i}| < 1$ 

$$\begin{split} \mathbb{E} \Big[ \sup_{g_{\theta} \in B(m',\mu)} \Big( \nu_n(g_{\theta}) - \nu_n^*(g_{\theta}) \Big)^2 \Big] \leq & \frac{1}{n} \left( \tau^{(2)}(q_n) \right)^2 \\ \leq & \frac{C_{\tau}^2}{n} \left( \frac{\log q_n}{q_n} \right)^{2\gamma - 2} \end{split}$$

where the last inequality follows from Proposition 3. Thus,

$$\begin{split} \mathbb{E}\Big[\sup_{g_{\theta}\in B(\widehat{m},\mu)} \Big(v_n(g_{\theta}) - v_n^*(g_{\theta})\Big)^2\Big] \leq & \sum_{m'\in\mathcal{M}_n} \mathbb{E}\Big[\sup_{g_{\theta}\in B(m',\mu)} \Big(v_n(g_{\theta}) - v_n^*(g_{\theta})\Big)^2\Big] \\ \leq & K_n \frac{C_{\tau}^2}{n} \left(\frac{\log q_n}{q_n}\right)^{2\gamma - 2} \\ \leq & \frac{A^2 C_{\tau}^2}{n^2}, \end{split}$$

using Assumption A5 and since  $K_n \leq n$ .

# 2/ Control of $\mathbb{E}[V_{\widehat{m}}^*]$ .

First, let us rewrite  $v_n^*(g_\theta)$  for  $g_\theta \in B(m', \mu)$ . Setting  $\mathbf{X}_t = (X_{t-1}^*, \dots, X_{t-D(S_{m'})}^*)^\top$ , we have



$$\begin{split} v_n^*(g_\theta) &= \frac{1}{n} \sum_{t=1}^n \xi_t g_\theta(\mathbf{X}_t^*) \\ &= \frac{1}{2 s_n q_n} \sum_{k=0}^{s_n - 1} \left( \sum_{i=1}^{q_n} \xi_{2kq_n + i} g_\theta(\mathbf{X}_{2kq_n + i}^*) + \sum_{i=1}^{q_n} \xi_{(2k+1)q_n + i} g_\theta(\mathbf{X}_{(2k+1)q_n + i}^*) \right) \\ &= v_{n,1}^*(g_\theta) + v_{n,2}^*(g_\theta) \end{split}$$

with

$$v_{n,1}^*(g_\theta) = \frac{1}{s_n} \sum_{k=0}^{s_n-1} v_{n,1,k}^*(g_\theta)$$
 and  $v_{n,2}^*(g_\theta) = \frac{1}{s_n} \sum_{k=0}^{s_n-1} v_{n,2,k}^*(g_\theta)$ 

where

$$\nu_{n,1,k}^*(g_\theta) = \frac{1}{2 \, q_n} \, \sum_{i=1}^{q_n} \, \xi_{2kq_n+i} \, g_\theta(\mathbf{X}_{2kq_n+i}^*) \quad \text{and} \quad \nu_{n,2,k}^*(g_\theta) = \frac{1}{2 \, q_n} \, \sum_{i=1}^{q_n} \, \xi_{(2k+1)q_n+i} \, g_\theta(\mathbf{X}_{(2k+1)q_n+i}^*).$$

Now let remark that  $v_{n,1}^*(g_\theta)$  and  $v_{n,2}^*(g_\theta)$  are both sum of  $s_n$  independent random variables by virtue of Proposition 4. Hence,

$$\begin{split} V_{\widehat{m}}^* & \leq \left(\sup_{g_{\theta} \in B(\widehat{m}, \mu)} 4 \left(v_{n, 1}^*(g_{\theta})\right)^2 - 4 x^2 n^{-1} D(S_{\widehat{m}})\right)_{+} \\ & + \left(\sup_{g_{\theta} \in B(\widehat{m}, \mu)} 4 \left(v_{n, 2}^*(g_{\theta})\right)^2 - 4 x^2 n^{-1} D(S_{\widehat{m}})\right)_{+}. \end{split}$$

As a consequence it is sufficient to study  $\mathbb{E}_1^* := \mathbb{E}\Big(\sup_{g \in B(\hat{m},\mu)} 4\left(v_{n,1}^*(g_\theta)\right)^2 - 4x^2 \, n^{-1} \, D(S_{\hat{m}})\Big)_+$  and the bound for  $\mathbb{E}\Big(\sup_{g_\theta \in B(\hat{m},\mu)} 4\left(v_{n,2}^*(g_\theta)\right)^2 - 4x^2 \, n^{-1} \, D(S_{\hat{m}})\Big)_+$  will follow by using analogous arguments.

# Bounding $\mathbb{E}_1^*$

Since the noise  $(\xi_t)$  is not bounded, the process  $v_{n,1}^*$  is not bounded either. Let's use the technique used in Comte and Genon-Catalot (2020) to overcome this difficulty. Therefore, we decompose  $\xi_t$  as

$$\xi_t = \eta_t + \epsilon_t, \quad \eta_t = \xi_t \Pi_{|\xi_t| \le k_n},$$

where  $k_n$  is a deterministic sequence or a constant to be chosen later. We then have

$$v_{n,1}^*(g_\theta) = v_{n,1}^*(g_\theta) + v_{n,2}^*(g_\theta), \text{ where}$$

$$v_{n,1}^*(g_\theta) = \frac{1}{s_n} \sum_{k=0}^{s_n-1} v_{n,1,k}^*(g_\theta) \quad \text{with} \quad v_{n,1,k}^*(g_\theta) = \frac{1}{2 q_n} \sum_{i=1}^{q_n} \eta_{2kq_n+i} g_\theta(\mathbf{X}_{2kq_n+i}^*) \quad \text{and} \quad v_{n,1,k}^*(g_\theta) = \frac{1}{2 q_n} \sum_{i=1}^{q_n} \eta_{2kq_n+i} g_\theta(\mathbf{X}_{2kq_n+i}^*)$$



$$v_{n,2}^*(g_\theta) = \frac{1}{s_n} \sum_{k=0}^{s_n-1} v_{n,2,k}^*(g_\theta) \quad \text{with} \quad v_{n,2,k}^*(g_\theta) = \frac{1}{2 q_n} \sum_{i=1}^{q_n} \epsilon_{2kq_n+i} g_\theta(\mathbf{X}_{2kq_n+i}^*).$$

Thus,

$$\mathbb{E}_{1}^{*} \leq 8 \,\mathbb{E}\left[\left(\sup_{g_{\theta} \in B(\widehat{m}, \mu)} \left(v_{n,1}^{*}(g_{\theta})\right)^{2} - 0.5 \,x^{2} \,n^{-1} \,D(S_{\widehat{m}})\right)_{+}\right] + 2 \,\mathbb{E}\left[\sup_{g_{\theta} \in B(\widehat{m}, \mu)} \left(v_{n,2}^{*}(g_{\theta})\right)^{2}\right]$$

$$\leq 8 \,\sum_{m' \in \mathcal{M}_{n}} \mathbb{E}\left[\left(\sup_{g_{\theta} \in B(m', \mu)} \left(v_{n,1}^{*}(g_{\theta})\right)^{2} - 0.5 \,x^{2} \,n^{-1} \,D(S_{m'})\right)_{+}\right]$$
(22)

$$+2\mathbb{E}\left[\sup_{g_{\theta}\in B(\widehat{m},\mu)}\left(v_{n,2}^{*}(g_{\theta})\right)^{2}\right]. \tag{23}$$

We start by bounding the term in (22). Let  $m' \in \mathcal{M}_n$ . In order to apply Theorem 2, one has to find M, H and v such that

$$\begin{split} \sup_{g_{\theta} \in B(m',\mu)} \left| v_{n,1,k}^*(g_{\theta}) \right| & \leq M, \quad \mathbb{E} \left[ \sup_{g_{\theta} \in B(m',\mu)} \left| v_{n,1}(g_{\theta}) \right|^2 \right] \leq H^2, \\ & \text{and } \sup_{g \in B(m',\mu)} \operatorname{Var} \left( v_{n,1,k}^*(g_{\theta}) \right) \right) \leq \nu. \end{split}$$

• Since the noise is bounded here and from the assumption **A1**, the process  $(X_t)$  is also bounded. Indeed, under **A1**, there exists  $(\phi_i^*)$  such that

$$X_t = \sum_{i=0}^{\infty} \phi_i^* \, \xi_{t-i} \quad \text{with} \quad \sum_{i=0}^{\infty} |\phi_i^*| < +\infty.$$

Therefore,  $|X_t| \le \Phi_0 k_n$  with  $\Phi_0 := \sum_{i=0}^{\infty} |\phi_i^*|$ . Moreover, for any  $g_\theta \in B(m', \mu)$ , we have

$$|g_{\theta}(\mathbf{X}_{t})| = |\sum_{i=1}^{D(S_{m'})} \theta_{i} X_{t-i}| \le \Phi_{0} k_{n} \sum_{i=1}^{D(S_{m'})} |\theta_{i}| < \Phi_{0} k_{n}.$$

As a result, we have

$$\sup_{g_{\theta} \in B(m',\mu)} |v_{n,1,k}^{*}(g_{\theta})| \leq \frac{1}{2 q_{n}} \sup_{g_{\theta} \in B(m',\mu)} \sum_{i=1}^{q_{n}} |\eta_{2kq_{n}+i} g_{\theta}(\mathbf{X}_{2kq_{n}+i}^{*})|$$

$$\leq \frac{\Phi_{0} k_{n}^{2}}{2} := M.$$

• Next, since the parameter set are compacts, there exists  $\theta_0 \in \Theta_m \cup \Theta_{m'}$  such that

$$\sup_{g_{\theta} \in B(m',\mu)} \left| v_{n,1}^*(g_{\theta}) \right|^2 = \left| v_{n,1}^*(g_{\theta_0}) \right|^2.$$

Moreover,



$$\begin{split} \mathbb{E} \big[ \big| v_{n,1}^*(g_{\theta_0}) \big|^2 \big] &= \frac{1}{s_n} \, \mathbb{E} \big[ \big| v_{n,1,0}^*(g_{\theta_0}) \big|^2 \big] \\ &= \frac{1}{4 \, s_n \, q_n^2} \, \sum_{i,j=1}^{q_n} \, \mathbb{E} \big[ \eta_i \, g_{\theta}(\mathbf{X}_i^*) \, \eta_j \, g_{\theta}(\mathbf{X}_j^*) \big] \\ &= \frac{1}{4 \, s_n \, q_n^2} \, \sum_{i=1}^{q_n} \, \mathbb{E} \big[ \big( \eta_i \, g_{\theta}(\mathbf{X}_i^*) \big)^2 \big] \\ &\leq \frac{\Phi_0^2 \, k_n^2}{2 \, n} \, \leq \, \frac{\Phi_0^2 \, k_n^2}{2 \, n} \, D(S_{m'}) \, := H^2 \end{split}$$

since  $D(S_{m'}) \ge 1$ .

• Lastly, as  $Var[X] \le E[X^2]$ , it follows from the previous series of equations

$$\operatorname{Var}\left(v_{n,1,0}^{*}(g_{\theta})\right)\right) \leq \mathbb{E}\left[\left|v_{n,1,0}^{*}(g_{\theta_{0}})\right|^{2}\right] \leq \frac{\Phi_{0}^{2} k_{n}^{2}}{4 q_{n}} := v.$$

As a consequence from Theorem 2 and taking  $\alpha = \frac{1}{2}(\frac{x^2}{2\Phi_0^2 k_+^2} - 1) > 0$ , we have

$$\begin{split} \mathbb{E}\left[\left(\sup_{g_{\theta}\in B(m',\mu)}\left(v_{n,1}^{*}(g_{\theta})\right)^{2}-0.5\,x^{2}\,n^{-1}\,D(S_{m'})\right)_{+}\right] \\ &\leq \frac{2}{K}\left(\frac{\Phi_{0}^{2}\,k_{n}^{2}}{4q_{n}}e^{-Kq_{n}D(S_{m'})(\frac{x^{2}}{2\Phi_{0}^{2}k_{n}^{2}}-1)}+\frac{49\Phi_{0}^{2}k_{n}^{4}}{4n^{2}KC^{2}(\alpha)}e^{-2\sqrt{2}KC(\alpha)\frac{\sqrt{n}\sqrt{D(S_{m'})}}{k_{n}}}\right). \end{split}$$

Hence, there exists a constant K' such that

$$\sum_{m' \in \mathcal{M}_n} \mathbb{E}\left[ \left( \sup_{g_{\theta} \in B(m', \mu)} \left( v_{n, 1}^*(g_{\theta}) \right)^2 - 0.5 \, x^2 \, n^{-1} \, D(S_{m'}) \right)_+ \right] \le \frac{K'}{n}. \tag{24}$$

• Now, let us upper bound the term in (23). For any  $m' \in \mathcal{M}_n$  and any  $g_{\theta} \in B(m', \mu)$ , we have

$$g_{\theta}(\mathbf{X}_{t}) = \sum_{i=1}^{D(S_{m'})} \theta_{i} X_{t-i} \leq \sup_{t-D(S_{m'}) \leq i < t} |X_{i}| \left(\sum_{i=1}^{D(S_{m'})} |\theta_{i}|\right) < \sup_{t-D(S_{m'}) \leq i < t} |X_{i}|.$$

Therefore,

$$v_{n,2,k}^{*}(g_{\theta}) = \frac{1}{2 q_{n}} \sum_{i=1}^{q_{n}} \xi_{2kq_{n}+i} g_{\theta}(\mathbf{X}_{2kq_{n}+i}^{*}) 
< \frac{1}{2 q_{n}} \sum_{i=1}^{q_{n}} |\xi_{2kq_{n}+i}| \sup_{2kq_{n}+i-K_{n} \le t < 2kq_{n}+i} |X_{t}^{*}| := Y_{k}^{*},$$
(25)

so that



$$\sup_{g \in B(\widehat{m}, \mu)} v_{n,2}^*(g_{\theta}) < \frac{1}{s_n} \sum_{k=0}^{s_n - 1} Y_k^*.$$

Let us notice that  $(Y_k^*)_k$  is a family of independent random variables as  $(v_{n,1,k}^*(g))_k$ . Thus, it follows

$$\begin{split} \mathbb{E} \Big[ \sup_{g_{\theta} \in B(\hat{m}, \mu)} \big| \nu_{n, 2}^{*}(g_{\theta}) \big|^{2} \Big] < & \frac{1}{s_{n}^{2}} \sum_{i, j = 0}^{s_{n} - 1} \mathbb{E} \big[ Y_{i}^{*} Y_{j}^{*} \big] \\ < & \frac{1}{s_{n}^{2}} \sum_{i = 0}^{s_{n} - 1} \mathbb{E} \big[ Y_{i}^{*2} \big] \\ = & \frac{1}{s_{n}} \mathbb{E} \big[ Y_{0}^{*2} \big]. \end{split}$$

Moreover,

$$\mathbb{E}\left[Y_{0}^{*2}\right] = \frac{1}{4 q_{n}^{2}} \sum_{i,j=1}^{q_{n}} \mathbb{E}\left[\left|\xi_{i}\right| \sup_{i-K_{n} \leq t < i} \left|X_{t}^{*}\right| \left|\xi_{j}\right| \sup_{j-K_{n} \leq t < j} \left|X_{t}^{*}\right|\right]$$

$$= \frac{1}{4 q_{n}^{2}} \sum_{i=1}^{q_{n}} \mathbb{E}\left[\left(\sup_{i-K_{n} \leq t < i} \left|X_{t}^{*}\right|\right)^{2}\right]$$

$$= \frac{\mu_{2}}{4 q_{n}},$$
(26)

where  $\mu_2 = \mathbb{E}[X_t^2] < \infty$ . It follows

$$\mathbb{E}\left[\sup_{g_{\theta} \in B(\hat{m}, \mu)} \left| v_{n,1}^*(g_{\theta}) \right|^2 \right] < \frac{\mu_2}{4 \, s_n \, q_n} = \frac{\mu_2}{2 \, n}. \tag{27}$$

Inequality (24) along with (27) yields to

$$\mathbb{E}_1^* \le \frac{8\,K'}{n} + \frac{\mu_2}{n}.$$

We conclude that there exists K > 0

$$\mathbb{E}[V_{\widehat{m}}] \le \frac{K}{n}.\tag{28}$$

Returning to (20), and taking expectation on both sides, it then follows

$$\left(1 - 2\frac{(1+y)}{x}\right) \mathbb{E}\left[\left\|f_{\hat{\theta}_{\hat{m}}}^{t} - f_{\theta^{*}}^{t}\right\|_{n}^{2}\right] \le \left(1 + 2\frac{(1+y^{-1})}{x}\right) \mathbb{E}\left[\left\|f_{\theta^{*}}^{t} - f_{\theta_{m}^{*}}^{t}\right\|_{n}^{2}\right] + 2\operatorname{pen}(S_{m}) + x\frac{K}{n}.$$
(29)

For 
$$y = \frac{x-2}{x+2} > 0$$
, so that  $1 + y = \frac{2x}{x+2}$  and  $1 + y^{-1} = \frac{2x}{x-2}$ , we obtain



$$\mathbb{E}\Big[ \| f_{\widehat{\theta}_{\widehat{m}}}^t - f_{\theta^*}^t \|_n^2 \Big] \le C(x) \left( \mathbb{E}\Big[ \| f_{\theta^*}^t - f_{\theta_m^*}^t \|_n^2 \Big] + 2 \operatorname{pen}(S_m) \right) + \frac{x(x+2)}{x-2} \frac{K}{n}$$
with  $C(x) = \frac{(x+2)^2}{(x-2)^2} > 1$ .

#### 4.2 Proof of Proposition 2

Since the collection  $\mathcal{M}_n$  is hierarchical, we have

$$\mathbb{P}(\Omega_n^c) \le \sum_{m \in \mathcal{M}_n} \mathbb{P}\left(\exists F_\theta \in S_m : \left| \frac{\|f_\theta\|_n^2}{\|f_\theta\|_\mu^2} - 1 \right| > \frac{1}{2}\right)$$

$$\le \sum_{m \in \mathcal{M}_n} \mathbb{P}(\Omega_m^c)$$

where

$$\Omega_m = \left\{ \left| \frac{\|F_\theta\|_n^2}{\|F_\theta\|_n^2} - 1 \right| \le \frac{1}{2} \quad \forall F_\theta \in S_m \right\}.$$

Let  $m \in \mathcal{M}_n$ . We have

$$\mathbb{P}(\Omega_m^c) \le \mathbb{P}\left(\sup_{F_\theta \in S_m} \left| \frac{\|F_\theta\|_n^2}{\|F_\theta\|_\mu^2} - 1 \right| > \frac{1}{2} \right).$$

Moreover,

$$\sup_{F_{\theta} \in S_m, \|F_{\theta}\|_{2}^{2} = 1} \left| \frac{\|F_{\theta}\|_{n}^{2}}{\|F_{\theta}\|_{u}^{2}} - 1 \right| > \frac{1}{2} \Longleftrightarrow \sup_{F_{\theta} \in S_m, \|F_{\theta}\|_{2}^{2} = 1} \left| \nu_{n}(F_{\theta}^{2}) \right| > \frac{1}{2}$$

with  $v_n(F_{\theta}^2) = n^{-1} \sum_{t=1}^n \left( (f_{\theta}^t)^2 - \mathbb{E}[(f_{\theta}^1)^2] \right)$ . Hence,

$$\mathbb{P} \left( \sup_{f \in S_m} \left| \frac{\|F_\theta\|_n^2}{\|F_\theta\|_n^2} - 1 \right| > \frac{1}{2} \right) \leq \mathbb{P} \left( \sup_{F_\theta \in S_m} \left| \nu_n(F_\theta^2) \right| > \frac{1}{2} \right).$$

For any  $F_{\theta} \in S_m$ , using the linearity we can write

$$\left(f_{\theta}(\mathbf{X}_{t})\right)^{2} = \left(\sum_{i=1}^{D_{m}} \theta_{i} X_{t-i}\right)^{2} = \sum_{i=1}^{D_{m}} \theta_{i} \theta_{j} X_{t-i} X_{t-j} = \theta^{\top} \widehat{\Sigma}_{m,t} \theta$$

where  $\theta = (\theta_1, \dots, \theta_{D_m})^{\mathsf{T}}$ ,  $\widehat{\Sigma}_{m,t} = Z_t^m (Z_t^m)^{\mathsf{T}}$  with  $Z_t^m = (X_{t-1}, \dots, X_{t-D_m})^{\mathsf{T}}$ . So that with  $\Sigma_m = \mathbb{E}[\widehat{\Sigma}_{m,t}]$ , it follows



$$\nu_n^*(F_\theta^2) = \frac{1}{n} \sum_{t=1}^n \theta^\top (\widehat{\Sigma}_t - \Sigma) \theta = \theta^\top (\widehat{\Sigma}_m - \Sigma_m) \theta,$$

where  $\hat{\Sigma} = n^{-1} \sum_{t=1}^{n} \hat{\Sigma}_{t}$ . As a result,

$$\sup_{F_{\theta} \in S_m, \|F_{\theta}\|_u^2 = 1} |v_n^*(F_{\theta}^2)| \le \|\widehat{\Sigma}_m - \Sigma_m\|_{\text{op}}.$$

Indeed,

$$\sup_{F_{\theta} \in S_m} |v_n^*(F_{\theta}^2)| = \sup_{\theta \colon \sum |\theta_i| < 1} \theta^{\mathsf{T}} (\widehat{\Sigma}_m - \Sigma_m) \theta = \sup_{\theta \colon \sum |\theta_i| < 1} \|\theta\|^2 \frac{\theta^{\mathsf{T}} (\widehat{\Sigma}_m - \Sigma_m) \theta}{\|\theta\|^2}$$

$$\leq \sup_{\theta \colon \|\theta\|^2 \le 1} \frac{\theta^{\mathsf{T}} (\widehat{\Sigma}_m - \Sigma_m) \theta}{\|\theta\|^2} = \|\widehat{\Sigma}_m - \Sigma_m\|_{\text{op}}$$

since  $-1 < \theta_i < 1$  ensures that  $\|\theta\|^2 \le \sum |\theta_i|$ . Hence,

$$\begin{split} \mathbb{P}\bigg(\sup_{F_{\theta} \in S_m, \|F_{\theta}\|_{\mu}^2 = 1} |\nu_n(F_{\theta}^2)| > \frac{1}{2}\bigg) \leq & \mathbb{P}\bigg(\|\widehat{\Sigma}_m - \Sigma_m\|_{\text{op}} > \frac{1}{2}\bigg) \\ \leq & \mathbb{P}\bigg(\|\widehat{\Sigma}_m - \widehat{\Sigma}_m^*\|_{\text{op}} > \frac{1}{4}\bigg) + \mathbb{P}\bigg(\|\widehat{\Sigma}_m^* - \Sigma_m\|_{\text{op}} > \frac{1}{4}\bigg) \\ = & : \mathbb{P}_1 + \mathbb{P}_2. \end{split}$$

Using Lemma 3 with u = 1/4 and by virtue of **A6**, it follows

$$\mathbb{P}_2 \le 2 \exp\left(-3\log n\right)$$

$$\le \frac{2}{n^3}.$$

Now let bound  $\mathbb{P}_1$ . We know that for a  $D_m \times D_m$  matrix A

$$||A||_{\text{op}} \le ||A||_{\infty} := \max_{1 \le i \le D_m} \sum_{i=1}^{D_m} |A_{ij}|.$$

Thus, from Markov's Inequality,



$$\begin{split} \mathbb{P}_{1} \leq & 4 \, \mathbb{E} \Big[ \big\| \widehat{\Sigma}_{m} - \widehat{\Sigma}_{m}^{*} \big\|_{\mathrm{op}} \Big] \\ \leq & 4 \, \mathbb{E} \Big[ \max_{1 \leq i \leq D_{m}} \sum_{j=1}^{D_{m}} \big| \big( \widehat{\Sigma}_{m} - \widehat{\Sigma}_{m}^{*} \big)_{i,j} \big| \Big] \\ \leq & 4 \, \sum_{j=1}^{D_{m}} \mathbb{E} \Big[ \big| \big( \widehat{\Sigma}_{m} - \widehat{\Sigma}_{m}^{*} \big)_{i_{0},j} \big| \Big] \\ \leq & 4 \, \sum_{j=1}^{D_{m}} \mathbb{E} \Big[ \big| X_{t-i_{0}} X_{t-j} - X_{t-i_{0}}^{*} X_{t-j}^{*} \big| \Big]. \end{split}$$

Moreover,  $|X_{t-i}X_{t-j} - X_{t-i}^*X_{t-j}^*| \le |X_{t-i}||X_{t-j} - X_{t-j}^*| + |X_{t-j}^*||X_{t-i} - X_{t-i}^*|$  so that with Cauchy–Schwarz's Inequality,

$$\begin{split} \mathbb{E}\left[\left|X_{t-i}X_{t-j} - X_{t-i}^*X_{t-j}^*\right|\right] &\leq 2 \left\|X_0\right\|_2 \left\|X_{t-1} - X_{t-1}^*\right\|_2 \\ &\leq 2 \left\|X_0\right\|_2 \tau^{(2)}(q_n). \end{split}$$

Hence, using Proposition 3, it follows

$$\begin{split} \mathbb{P}_{1} \leq & 8 \, \big\| X_{0} \big\|_{2} D_{m} \, \tau^{(2)}(q_{n}) \\ \leq & 8 \, \big\| X_{0} \big\|_{2} D_{m} \, \, C_{\tau} \left( \frac{\log q_{n}}{q_{n}} \right)^{\gamma - 1}. \end{split}$$

Moreover, since  $\gamma \geq 8$  and from assumption **A5**, one can find some constant A' such that

$$\left(\frac{\log q_n}{q_n}\right)^{\gamma-1} \le \frac{A'}{n^4}.$$

As a result, with  $c_0 := 8 \|X_0\|_2 C_{\tau} A'$ , it holds

$$\mathbb{P}_1 \le \frac{c_0}{n^3}.$$

As a consequence,

$$\mathbb{P}(\Omega_n^c) \leq \frac{2+c_0}{n^3}.$$

## 4.3 Proof of Proposition 1

**Proof** The proof of the will be based on the relation between the spectral density function and the maximum eigenvalues of the variance covariance matrix.

Denote by  $u \in \mathbb{R}^{D_m}$  the normalized eigenvector associated with the largest eigenvalue  $\lambda_{\max}(\Sigma_m)$ . Hence,



$$\lambda_{\max}(\Sigma_m) = u^{\mathsf{T}} \Sigma_m u = \sum_{j,k=1}^{D_m} u_j \, r(j-k) \, u_k = \int_{-\pi}^{\pi} g(\lambda) \sum_{j,k=1}^{D_m} u_j \, e^{i(j-k)\lambda} \, u_k d\lambda$$

$$= \int_{-\pi}^{\pi} g(\lambda) \left| \sum_{j=1}^{D_m} u_j \, e^{ij\lambda} \right|^2 d\lambda \leq \sup_{-\pi \leq \lambda < \pi} g(\lambda) \int_{-\pi}^{\pi} \left| \sum_{j=1}^{D_m} u_j \, e^{ij\lambda} \right|^2 d\lambda$$

$$\leq \sup_{-\pi \leq \lambda < \pi} g(\lambda),$$

since, using Parseval identity,  $\int_{-\pi}^{\pi} \left| \sum_{j=1}^{D_m} u_j e^{ij\lambda} \right|^2 d\lambda = \sum_{j=1}^{D_m} u_j^2 = 1$ . But, from Lemma 2 and since  $\gamma \ge 2$ , it follows

$$\left| \sup_{-\pi \le \lambda < \pi} g(\lambda) \right| \le \frac{1}{2\pi} \sum_{h \in \mathbb{Z}} |r(h)|$$
$$\le \frac{C}{\pi} \sum_{h=0}^{+\infty} \frac{1}{(h+1)^{\gamma}} < \infty.$$

Given that  $\Sigma_m$  is symmetric, it follows

$$\left\| \Sigma_m \right\|_{\text{op}} = \lambda_{\max}(\Sigma_m) \le \frac{C}{\pi} \sum_{k=0}^{+\infty} \frac{1}{(h+1)^{\gamma}},$$

which concludes the proof of (11).

Now we end by the proof of (12). Reasoning as above, and by virtue of A4, one can show that

$$\lambda_{\min}(\Sigma_m) \ge \inf_{-\pi < \lambda < \pi} g(\lambda) \ge a$$

which yields to

$$\left\|\Sigma_m^{-1}\right\|_{\text{op}} = \frac{1}{\lambda_{\min}(\Sigma_m)} \le \frac{1}{a},$$

so that (12) is established.

#### 4.4 Technical lemmas

**Lemma 1** Assume **A1** holds and  $(X_t)$  the mixing stationary solution of (1). Then, the process  $(X_t)$  is mixing and

$$\tau_{\mathbf{X},\infty}^{(1)}(r) \le K_n \, \tau_{X,\infty}^{(1)}(r-1). \tag{30}$$



**Proof** Let set by  $\mathcal{M}_{\mathbf{X}}^i = \sigma(\mathbf{X}_t, t \leq i)$  and  $\mathcal{M}_X^i = \sigma(X_t, t \leq i)$  for an integer i. One would like to bound  $\tau^{(1)}(\mathcal{M}_{\mathbf{X}}^i, (\mathbf{X}_{i_1}, \dots, \mathbf{X}_{i_r}))$  for  $j_k > \dots > j_1 \geq i + r$ .

would like to bound  $\tau^{(1)}\big(\mathcal{M}_{\mathbf{X}}^i, \big(\mathbf{X}_{j_1}, \dots, \mathbf{X}_{j_k}\big)\big)$  for  $j_k > \dots > j_1 \geq i + r$ . Let assume that the universe  $\Omega$  is rich enough so that, one can find  $\mathbf{X}_{j_l}^* = \big(X_{j_l-1}^*, \dots, X_{j_l-K_n}^*\big)^{\mathsf{T}}$  with  $l=1,\dots,k$  verifying

- 1.  $(\mathbf{X}_{i_1}^*, \dots, \mathbf{X}_{i_k}^*)$  is distributed as  $(\mathbf{X}_{j_1}, \dots, \mathbf{X}_{j_k})$  and independent of  $\mathcal{M}_{\mathbf{X}}^i$ ;
- 2.  $(X_{j_1-1}^*, \dots, X_{j_k-1}^*)^{\mathsf{T}}$  is distributed as  $(X_{j_1-1}, \dots, X_{j_k-1})^{\mathsf{T}}$  and independent of  $\mathcal{M}_X^i$ .

As a result,

$$\tau^{(1)}(\mathcal{M}_{\mathbf{X}}^{i}, (\mathbf{X}_{j_{1}}, \dots, \mathbf{X}_{j_{k}})) \leq \sum_{l=1}^{k} \|\mathbf{X}_{j_{l}} - \mathbf{X}_{j_{l}}^{*}\|_{1} = \sum_{l=1}^{k} \sum_{t=1}^{K_{n}} \mathbb{E}[|X_{j_{t-t}} - X_{j_{t-t}}^{*}|]$$

$$\leq K_{n} \sum_{l=1}^{k} \mathbb{E}[|X_{j_{t-1}} - X_{j_{t-1}}^{*}|]$$

$$= K_{n} \|(X_{j_{1}-1}, \dots, X_{j_{k}-1})^{\top} - (X_{j_{1}-1}^{*}, \dots, X_{j_{k}-1}^{*})^{\top}\|_{1}$$

$$= K_{n} \tau^{(1)}(\mathcal{M}_{X}^{i}, (X_{j_{1}-1}, \dots, X_{j_{k}-1})).$$

This fact along with the definition of  $\tau_{\mathbf{X},\infty}^{(1)}(r)$  leads to (30).

**Lemma 2** Under **A1** with  $|\theta_{t}^{*}| = O(t^{-\gamma})$  where  $\gamma > 1$ , we have

$$r(h) = \mathbb{E}[X_0 X_h] = O((h+1)^{-\gamma})$$

**Proof** By virtue of **A1**, the process  $(X_t)_t$  is causal; that is there exists  $(\phi_i)_{i\in\mathbb{N}}$  such that  $X_t = \sum_{i=0}^{+\infty} \phi_i \, \xi_{t-i}$  with  $\sum_{i=0}^{+\infty} |\phi_i| < \infty$ . The sequence  $(\phi_i)_{i\in\mathbb{N}}$  is given by the relation  $\phi(z) = \sum_{i=0}^{+\infty} \phi_i \, z^i = \frac{1}{\theta(z)}$  with  $\theta(z) = 1 - \sum_{i=0}^{+\infty} \theta_i^* \, z^i$ . Equating coefficients of  $z_i, j = 0, 1, \ldots$ , we find that  $\phi_0 = 1$  and for  $i \ge 1$ 

$$\phi_i = \sum_{j=1}^i \theta_j^* \, \phi_{i-j}.$$

This fact allows us to deduce that the sequences  $(\phi_i)_{i \in \mathbb{N}}$  and  $(\theta_i^*)_{i \in \mathbb{N}}$  decay at the same rate. Therefore, since  $|\theta_t^*| = O((t+1)^{-\gamma})$ , there exists  $h_0 \in \mathbb{Z}$  such that for any  $h \ge h_0$ , it holds  $|\phi_t| \le C(t+1)^{-\gamma}$  for some constant C > 0. Thus,



$$\begin{split} r(h) &= \sum_{j=0}^{\infty} \phi_j \, \phi_{j+h} \\ &\leq C^2 \sum_{j=0}^{\infty} \frac{1}{(j+1)^{\gamma}} \, \frac{1}{(j+h+1)^{\gamma}} \\ &\leq C^2 \, (h+1)^{-\gamma} \, \sum_{j=0}^{\infty} \frac{1}{(j+1)^{\gamma}} &\leq C^2 \frac{\pi^2}{6} \, (h+1)^{-\gamma}, \end{split}$$

where the last inequality follows from the fact that  $\gamma \ge 2$  and that established the Lemma.

**Lemma 3** Under assumptions **A2**, it holds for any model  $m \in \mathcal{M}_n$ , and for all u > 0

$$\mathbb{P}\Big(\|\widehat{\Sigma}_m^* - \Sigma_m\|_{\text{op}} \ge u\Big) \le 2 \exp\left\{-\frac{s_n}{2} \min\left\{\left(\frac{u}{16D_m \sigma_0^2}\right)^2, \frac{u}{32D_m \sigma_0^2}\right\}\right\}$$

**Proof** One can write for a matrix A

$$||A||_{\text{op}} = \max_{v \cdot ||v||=1} |v^{\mathsf{T}} A v| = |v_0^{\mathsf{T}} A v_0|.$$

Therefore, one can find a vector  $v_0 \in \mathbb{R}^{D_m}$  with  $||v_0|| = 1$  such that

$$\mathbb{P}\Big(\|\widehat{\Sigma}_m^* - \Sigma_m\|_{\text{op}} \ge u\Big) = \mathbb{P}\Big(|v_0^\top (\widehat{\Sigma}_m^* - \Sigma_m) v_0| \ge u\Big).$$

But,

$$\begin{aligned} v_0^{\top} \big( \widehat{\Sigma}_m^* - \Sigma_m \big) v_0 &= \frac{1}{n} \sum_{t=1}^n \big( v_0^{\top} \, \widehat{\Sigma}_{m,t}^* v_0 - v_0^{\top} \Sigma_m v_0 \big) \\ &= \frac{1}{n} \sum_{t=1}^n \big( v_0^{\top} \, (Z_t^{*m}) \, (Z_t^{*m})^{\top} v_0 - v_0^{\top} \Sigma_m v_0 \big) \\ &= \frac{1}{n} \sum_{t=1}^n \big( Y_t^2 - \mathbb{E}[Y_t^2] \big) \end{aligned}$$

with  $Y_t = v_0^{\mathsf{T}} Z_t^m = \sum_{i=1}^{D_m} v_0^i X_{t-i}^*$ . From **A2**,  $Y_t$  is  $\mathrm{SG}(D_m \, \sigma_0^2)$ . Therefore,  $Y_t^2$  is  $\mathrm{SE}(256 \, D_m^2 \, \sigma_0^4, 16 \, D_m \sigma_0^2)$  (where SE stands for Sub-Gaussian and SE for Sub-Exponential).

Moreover, we can write



$$\begin{split} v_0^{\mathsf{T}} \left( \hat{\Sigma}_m^* - \Sigma_m \right) v_0 &= \frac{1}{n} \sum_{t=1}^n \left( Y_t^2 - \mathbb{E}[Y_t^2] \right) \\ &= \frac{1}{s_n} \sum_{k=0}^{s_n - 1} \left( \frac{1}{2q_n} \sum_{i=1}^{q_n} \left( Y_{2kq_n + i}^2 - \mathbb{E}[Y_1^2] \right) \right) \\ &+ \frac{1}{s_n} \sum_{k=0}^{s_n - 1} \left( \frac{1}{2q_n} \sum_{i=1}^{q_n} \left( Y_{(2k+1)q_n + i}^2 - \mathbb{E}[Y_1^2] \right) \right) \\ &= \mathbf{Y}_1 + \mathbf{Y}_2. \end{split}$$

Therefore,

$$\mathbf{Y}_{1} = \frac{1}{s_{n}} \sum_{k=0}^{s_{n}-1} \mathbf{Y}_{1,k} \quad \text{and} \quad \mathbf{Y}_{2} = \frac{1}{s_{n}} \sum_{k=0}^{s_{n}-1} \mathbf{Y}_{2,k} \quad \text{with}$$

$$\mathbf{Y}_{1,k} = \frac{1}{2q_{n}} \sum_{i=1}^{q_{n}} \left( Y_{2kq_{n}+i}^{2} - \mathbb{E}[Y_{1}^{2}] \right) \quad \text{and} \quad \mathbf{Y}_{2,k} = \frac{1}{2q_{n}} \sum_{i=1}^{q_{n}} \left( Y_{(2k+1)q_{n}+i}^{2} - \mathbb{E}[Y_{1}^{2}] \right).$$

 $\{\mathbf{Y}_{1,k}\}$  and  $\{\mathbf{Y}_{2,k}\}$  are independent random vectors by virtue of Proposition 4. Now, let us show that  $\mathbf{Y}_{i,k}$  are sub-exponentials. For  $\lambda$  such that  $|\lambda| < \frac{1}{16D_m\sigma_0^2}$ , and denoting  $w_i = Y_{2ka+i}^2 - \mathbb{E}[Y_1^2]$ , we have

$$\begin{split} \mathbb{E}\left[e^{\lambda \mathbf{Y}_{1,k}}\right] &= \mathbb{E}\left[\exp\left(\frac{1}{2q_n}\sum_{i=1}^{q_n}\lambda\,w_i\right)\right] \\ &= \mathbb{E}\left[\Pi_{i=1}^{q_n}\exp\left(\frac{\lambda\,w_i}{2q_n}\right)\right] \\ &= \mathbb{E}\left[\Pi_{i=1}^{q_n}\left(\exp\left(\frac{\lambda\,w_i}{2}\right)\right)^{1/q_n}\right] \\ &\leq \Pi_{i=1}^{q_n}\left(\mathbb{E}\left[\exp\left(\frac{\lambda\,w_i}{2}\right)\right]\right)^{1/q_n} \\ &\leq e^{\frac{\lambda^2}{2}64D_m^2\,\sigma_0^4}, \end{split}$$

where we have used Hölder's Inequality. Hence,  $\mathbf{Y}_{1,k}$  is  $SE(64 D_m^2 \sigma_0^4, 16 D_m \sigma_0^2)$ . As a result, using exponential inequalities for SE random variables, it follows

$$\mathbb{P}\left(\mathbf{Y}_1 \ge u/2\right) \le \exp\left\{-\frac{s_n}{2} \min\left\{\left(\frac{u}{16D_m \sigma_0^2}\right)^2, \frac{u}{32D_m \sigma_0^2}\right\}\right\}$$

so that

$$\mathbb{P}\left(\left|v_0^{\top}\left(\widehat{\Sigma}_m^* - \Sigma_m\right)v_0\right| \ge u/2\right) \le 2 \exp\left\{-\frac{s_n}{2} \min\left\{\left(\frac{u}{16D_m \sigma_0^2}\right)^2, \frac{u}{32D_m \sigma_0^2}\right\}\right\}.$$



**Lemma 4** Assume **A3** holds, then  $\hat{\Sigma}_m$  is a.e. invertible. Also,  $\Sigma_m$  is invertible.

**Proof** We can write  $\widehat{\Sigma}_m = \mathbf{M}_m^{\mathsf{T}} \mathbf{M}_m$  with  $\mathbf{M}_m = \left[ X_{i-1}, \dots, X_{i-D_m} \right]_{i=1}^n$ . By virtue of **A3**,  $\mathbf{M}_m$  is of full rank which implies the a.e. invertibility of  $\widehat{\Sigma}_m$ .

Moreover,  $\Sigma_m = \mathbb{E}\left[\widehat{\Sigma}_m\right] = \mathbb{E}\left[Z_0^m(Z_0^m)^\top\right]$  with  $Z_0^m = (X_{-1}, \dots, X_{-D_m})^\top$ . Let  $\mathbf{u} \in \mathbb{R}^{D_m}$ , it follows  $\mathbf{u}^\top \Sigma_m \mathbf{u} = \mathbb{E}\left[((Z_0^m)^\top \mathbf{u})^2\right] \geq 0$ . Let show that whenever the equality holds  $(\mathbf{u}^\top \Sigma_m = 0), \mathbf{u} = 0$ .

Since  $((Z_0^m)^\mathsf{T} \mathbf{u})^2 \ge 0$ , its expectation vanishes if and only if  $(Z_0^m)^\mathsf{T} \mathbf{u} = 0$  a.e. which yields to  $\mathbf{u} = 0$  by **A3**. Hence,  $\Sigma_m$  is positive definite and then invertible.

## 5 Theoretical tools

The following Proposition that is a consequence of Theorem 3.1 in Doukhan and Wintenberger (2008) gives a link between the  $\tau$ -mixing coefficients of the process  $(X_t)_{t \in \mathbb{Z}}$  and the coefficients  $\theta_i^*$  of model (3).

**Proposition 3** Assume **A1** holds and if  $|\theta_t^*| = O(t^{-\gamma})$  with  $\gamma > 1$ , there exists a  $\tau$ -weakly dependent stationary solution of (1) and a constant  $C_{\tau} > 0$  such that for r > 0

$$\tau_{X,\infty}^{(2)}(r) \le C_{\tau} \left(\frac{\log r}{r}\right)^{\gamma - 1}.\tag{31}$$

**Proof** With  $G(x, \xi_0) = \sigma \xi_0 + f_{\theta^*}(x)$  for any  $x \in \mathbb{R}^{\infty}$ , it holds

$$\left\|G(x,\xi_0) - G(y,\xi_0)\right\|_2 = \left|f_{\theta^*}(x) - f_{\theta^*}(y)\right| \le \sum_{i=1}^{\infty} |\theta_i^*| |x_i - y_i|.$$

Therefore, (31) is a straightforward application of Theorem 3.1 in Doukhan and Wintenberger (2008).

The following proposition allows us to obtain the block independence property.

**Proposition 4** Let  $(X_t)_{t \in \mathbb{Z}}$  be the stationary mixing process obtained in Proposition 3. Let also  $s_n, q_n, A_k, B_k$  defined as above for  $k = 0, ..., s_n - 1$ . There exist random vectors  $A_k^* = (\mathbf{X}_{2kq_n+1}^*, ..., \mathbf{X}_{(2k+1)q_n}^*), B_k^* = (\mathbf{X}_{(2k+1)q_n+1}^*, ..., \mathbf{X}_{(2k+2)q_n}^*)$  such that:

- 1. For  $k = 0, ..., s_n 1$ ,  $A_k^*$  has the same law as  $A_k$ , also  $B_k^*$  and  $B_k$ .
- 2. The random vectors  $(A_k^*)_{0 \le k \le s_n-1}$  are independent and so are the vectors  $(B_k^*)_{0 \le k \le s_n-1}$ .



3. 
$$\|A_k - A_k^*\|_1 \le q_n K_n \tau_{X,\infty}^{(1)}(q_n)$$
 and  $\|B_k - B_k^*\|_1 \le q_n K_n \tau_{X,\infty}^{(1)}(q_n).$ 

The next Theorem is a Talagrand's Inequality given in Klein and Rio (2005).

**Theorem 2** Let  $Y_1, ..., Y_n$  be independent random variables and let  $\mathcal{F}$  be a countable class of uniformly bounded measurable functions. Then, for all  $\alpha > 0$ ,

$$\mathbb{E}\left[\sup_{g \in \mathcal{F}} |\eta_n(g)|^2 - 2(1 + 2\alpha)H^2\right]_{+} \leq \frac{2}{K} \left(\frac{v}{n} e^{-K\alpha \frac{nH^2}{v}} + \frac{49M^2}{4Kn^2C^2(\alpha)} e^{-\frac{2\sqrt{2}KC(\alpha)\sqrt{\alpha}}{7\sqrt{2}} \frac{nH}{M}}\right)$$

with  $\eta_n(g) = n^{-1} \sum_{t=1}^n (g(Y_t) - \mathbb{E}[g(Y_t)])$  for any  $g \in \mathcal{F}$ ;

$$C(\alpha) = (\sqrt{1+\alpha} - 1) \wedge 1, K = 1/6$$

$$\sup_{g \in \mathcal{F}} \|g\|_{\infty} \le M, \quad \mathbb{E} \Big[ \sup_{g \in \mathcal{F}} |\eta_n(g)| \Big] \le H, \quad \sup_{g \in \mathcal{F}} \frac{1}{n} \sum_{t=1}^n \operatorname{Var}(g(Y_t)) \le v.$$

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