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# SEMIPARAMETRIC INFERENCE IN CORRELATED LONG MEMORY SIGNAL PLUS NOISE MODELS 

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#### Abstract

This paper proposes an extension of the log periodogram regression in perturbed long memory series that accounts for the added noise, also allowing for correlation between signal and noise, which represents a common situation in many economic and financial series. Consistency (for $d<1$ ) and asymptotic normality (for $d<3 / 4$ ) are shown with the same bandwidth restriction as required for the original $\log$ periodogram regression in a fully observable series, with the corresponding gain in asymptotic efficiency and faster convergence over competitors. Local Wald, Lagrange Multiplier and Hausman type tests of the hypothesis of no correlation between the latent signal and noise are also proposed.


JEL classification: C22; C13
Keywords: long memory; signal plus noise; log periodogram regression; semiparametric inference.

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## 1 INTRODUCTION

The analysis of economic and financial time series often needs to deal with situations in which a variable is not directly observable because it suffers contamination from some noise. The variable of interest in this case is a latent signal and the contaminating noise is usually incorporated in an additive form (perhaps after a logarithmic transformation) such that the observed series is of the form

$$
\begin{equation*}
z_{t}=y_{t}+u_{t} \tag{1}
\end{equation*}
$$

where $y_{t}$ is the latent signal and $u_{t}$ is the perturbing noise, usually considered to be weak dependent or even white noise. It is assumed without loss of generality that $z_{t}$ has a zero mean since the results described hereafter would not be altered by the addition of a non zero constant.

This situation is quite common in economic time series where the latent variable often shows strong persistence (Granger, 1966). For example economic mechanisms where the short run and long run behaviour of the series are affected by different factors may give rise to a series such as (1) where $y_{t}$ and $u_{t}$ represent the short and long run effects respectively. A second example is the measurement error which is concomitant to many economic variables and makes the variable of interest into a latent signal. A similar situation arises also in rational expectation models where the ex ante variable $y_{t}$ may exhibit long range dependence (e.g. Sun and Phillips, 2004, for the analysis of the long run Fisher equation).

In financial time series the noise may emerge in stock price series as a result of price discreteness or microstructure effects. In this case the noise represents the discrepancy between transaction prices and implicit efficient prices (Hasbrouck, 1993) and causes the weak autocorrelation empirically observed in many financial return series. This weak linear dependence of the returns is far from strong persistent. However, the noise in prices brings about a noisy series of squared returns, often used as proxies of the volatility, that masks the strong persistence in the (latent) volatility dynamics (Andersen and Bollerslev, 1998). The realized volatility (RV), built upon a summation of high frequency squared returns over a specified period, inherits the effects of this microstructure noise and can be considered to conform the specification in (1) with $y_{t}$ the latent volatility and $u_{t}$ emerging because of the microstructure noise. In fact a strong persistence has been often found in RV series,
so that ARFIMA models are usually employed (see Andersen et al, 2003, Deo et al, 2006, or Lieberman and Phillips, 2008). Alternatively, the strong persistence empirically found in the volatility of financial series can be modeled by means of stochastic volatility models. The Long Memory in Stochastic Volatility (LMSV) introduced independently by Harvey (1998) and Breidt et al. (1998) characterizes the returns as $r_{t}=\exp \left(y_{t} / 2\right) \varepsilon_{t}$ with $y_{t}$ the long memory component and $\varepsilon_{t}$ an independent white noise. The logs of the squared returns are of the form in (1) with a noise $u_{t}=\log \varepsilon_{t}^{2}$. In this context long and short-lived shocks are jointly considered in $y_{t}$. However, different types of news may affect volatility in different ways such that short and long run effects could be separately incorporated in the volatility specification as $r_{t}=\exp \left(\left(a_{1} y_{t}+a_{2} g_{t}\right) / 2\right) \varepsilon_{t}$ for some weak dependent $g_{t}$ and constants $a_{1}$ and $a_{2}$ (Bollerslev and Jubinski, 1999, Veiga 2006). In this case $\log r_{t}^{2}$ is the sum of a long memory process and a weak dependent noise $g_{t}+u_{t}$, rather than a white noise perturbing variable. Other log ARCH type models are also of this form with a latent signal corresponding to the persistent underlying volatility component.

Independence between signal and noise is usually assumed in the analysis of these series. However, this hypothesis is in many cases hard to sustain. The factors that affect the short run behaviour of a series might have also some effect on its long run behaviour and viceversa. Additionally, the potential correlation between a true economic variable and a measurement error has been demonstrated in a number of papers. For example, Bound et al. (1994) mention that the assumption of uncorrelation between the latent variable and the measurement error in labour market data "reflects convenience rather than conviction". Similarly, De Jong et al. (1998) show that transaction costs and lagged adjustment to information give rise to correlation between the underlying price and noise in stock price series that may transfer to the RV such that correlation between noise and latent volatility is expected in RV series. If the volatility is considered by means of stochastic volatility models, the leverage effect typically found in financial time series may introduce correlation between the latent volatility and the added noise. Additionally, if short and long-lived shocks to volatility are modeled separately their independence is at least debatable because there may exist shocks with both short and long run effects.

Independence of signal and noise is at least a debatable issue and the specification in (1) with correlated $y_{t}$ and $u_{t}$ can be considered as appropriate for many economic and financial
series. Ignoring either the presence of the noise or the correlation between signal and noise in the estimation procedure may mask the strong persistence of the latent signal. Considering both the noise and the correlation, the spectral density of $z_{t}$ in (1) can be expressed in terms of the spectral and cross-spectral densities of $y_{t}$ and $u_{t}$ such that

$$
\begin{equation*}
f_{z}(\lambda)=f_{y}(\lambda)+f_{u}(\lambda)+2 \operatorname{Re} f_{y u}(\lambda) \tag{2}
\end{equation*}
$$

where $\operatorname{Re}(a)$ denotes the real part of $a$. The cross spectral density only arises in the case of non null correlation of signal and noise.

The long memory of $y_{t}$ determines the behaviour of the spectral density at the origin such that $f_{y}(\lambda) \sim C_{y} \lambda^{-2 d}$ as $\lambda \rightarrow 0$ for a positive constant $C_{y}$. If the added noise does not show persistence or has less memory than $y_{t}$, and under a reasonable correlation structure between $y_{t}$ and $u_{t}$ (see the assumptions below), the long memory property of $y_{t}$ transmits to $z_{t}$ and the spectral density of the observable $z_{t}$ shares the divergency of $f_{y}(\lambda)$ at the origin with the same memory parameter $d>0$. This spectral property entitles the estimation of the memory parameter of the latent signal using semiparametric or local techniques originally proposed for fully observable long memory series, which only consider spectral behaviour around frequency zero. However, the added noise affects the properties of these estimators, inducing a large bias which limits the efficiency by compelling the use of frequencies very close to the origin. This effect has been analyzed by Deo and Hurvich (2001) and Arteche (2004) for the log-periodogram regression and the local Whittle estimators respectively. To reduce this bias Sun and Phillips (2003), Hurvich et al. (2005) and Arteche (2006) propose modifications of both estimators that include the added noise in the estimation procedures. Sun and Phillips (2003) and Arteche (2006) consider only the case of independent signal and noise in a log periodogram regression context. Hurvich et al. (2005) extend the local Whittle estimator by including explicitly in the estimation procedure both the added white noise and the potential correlation between signal and noise by incorporating terms that account for the spectral density of $u_{t}$ and the non null cross spectral density of $y_{t}$ and $u_{t}$. However, their proposal does not account fully for the correlation of signal and noise, which limits the asymptotic efficiency and rate of convergence further than claimed (see the corrigendum Hurvich et al., 2008).

Here we extend the log periodogram regression in analogous directions but account for more general correlation structures. As in other similar extensions, the local Whittle type
estimators dominate those based on a log periodogram regression in an asymptotic mean squared error sense. However, the local Whittle extension of Hurvich et al. (2005) is here improved in three directions. First, we allow for a more general specification of the correlation of signal and noise, covering more realistic situations. Second, the possibility of non contemporaneous correlation is also considered. Finally, weak dependence of the noise is allowed, which is particularly relevant not only in economic series but even in some extensions of the LMSV such as those in Bollerslev and Jubinski (1999) and Veiga (2006).

We also consider the nonstationary case as in Velasco (1999) and Hurvich et al. (2005). For that purpose we define the nonstationary $y_{t}$ as $y_{t}=y_{0}+\sum_{t=1}^{t} v_{t}$ where $y_{0}$ is a random variable not depending on $t$ and $v_{t}$ is weakly stationary with memory parameter $d-1 \in$ $[-1 / 2,0)$. The pseudo spectral density function of $y_{t}$ is then

$$
f_{y}(\lambda)=|1-\exp (i \lambda)|^{-2} f_{v}(\lambda) \sim C_{v} \lambda^{-2 d} \text { as } \lambda \rightarrow 0
$$

for $d \in[1 / 2,1)$. The pseudo spectral density function of $z_{t}$ can then be written in terms of the spectral density function of $v_{t}$ and the cross spectral density of $v_{t}$ and $u_{t}$ as

$$
f_{z}(\lambda)=|1-\exp (i \lambda)|^{-2} f_{v}(\lambda)+f_{u}(\lambda)+2 \operatorname{Re}\left\{(1-\exp (i \lambda))^{-1} f_{v u}(\lambda)\right\}
$$

or in terms of pseudo spectral and cross spectral densities as in $(2)$ with $f_{y u}(\lambda)=(1-$ $\exp (i \lambda))^{-1} f_{v u}(\lambda)$.

The rest of the paper is organized as follows. Section 2 describes the assumptions needed in our analysis and illustrates them with some primary examples. Section 3 introduces the proposed estimator and Section 4 shows its asymptotic properties, particularly consistency and asymptotic normality. Section 5 proposes Wald, Lagrange Multiplier and Hausman type tests of the hypothesis of no correlation between signal and noise. The finite sample performance of the proposed estimators and testing procedures are examined in Section 6 and in Section 7 they are applied to a series of daily RV of S\&P500 futures index. The technical details are relegated to the Appendices.

## 2 BASIC ASSUMPTIONS AND EXAMPLES

Assumption 1: The signal $y_{t}$ is a Gaussian process with a (pseudo) spectral density
satisfying
a) $f_{y}(\lambda)=C_{y} \lambda^{-2 d_{0}}\left(1+O\left(\lambda^{\beta_{1}}\right)\right)$
b) $f_{y}(\lambda)=C_{y} \lambda^{-2 d_{0}}\left(1+G_{y} \lambda^{\beta_{1}}+O\left(\lambda^{\beta_{1}+\iota}\right)\right)$
as $\lambda \rightarrow 0^{+}$, for some $\iota>0$, finite positive $C_{y}$, finite $G_{y}, 0<d_{0}<1, \beta_{1}>1+d_{0}$ and in a neighbourhood of the origin $f_{y}(\lambda)$ is differentiable with first derivative $O\left(\lambda^{-1-2 d_{0}}\right)$.

Assumption 2: The added noise $u_{t}$ is Gaussian with a spectral density satisfying

$$
f_{u}(\lambda)=f_{u}(0)\left(1+O\left(\lambda^{\beta_{2}}\right)\right)
$$

as $\lambda \rightarrow 0^{+}, \beta_{2}>\beta_{1}-2 d_{0}$, for a positive finite $f_{u}(0)$ and in a neighbourhood of the origin $f_{u}(\lambda)$ is differentiable with first derivative $O\left(\lambda^{-1}\right)$.

Assumption 3: As $\lambda \rightarrow 0^{+}$, the (pseudo) cross spectral density function of $y_{t}$ and $u_{t}$ satisfies

$$
\operatorname{Re}\left(f_{y u}(\lambda)\right)=\lambda^{-d_{0}}\left(C_{y u} \cos \left[d_{0}\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)\right]+G_{y u} \lambda \sin \left[d_{0}\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)\right]+O\left(\lambda^{\beta_{3}}\right)\right)
$$

for finite constants $C_{y u}$ and $G_{y u}, \beta_{3}>\beta_{1}-d_{0}$ and in a neighbourhood of the origin $f_{y u}(\lambda)$ is differentiable with first derivative $O\left(\lambda^{-1-d_{0}}\right)$.

Assumption 1 imposes a particular spectral behaviour of $y_{t}$ around zero, slightly relaxing Assumption 2 in Sun and Phillips (2003) and allowing for a non stationary $y_{t}$. The local specification in a) is required for consistency and b) is needed for the asymptotic normality. As in Henry and Robinson (1996) this local specification permits us to obtain the leading part of the asymptotic bias of local estimators of $d_{0}$ in terms of $G_{y}$. Only positive values of $d_{0}$ are considered. The condition $d_{0}>0$ guarantees that the long memory of $y_{t}$ is transferred to $z_{t}$ since by Assumption $2 u_{t}$ is weak dependent (e.g. stationary ARMA), which in view of the economic and financial examples given in the introduction is the empirically most interesting case. For $d_{0} \leq 0$ the persistence of $z_{t}$ would be that of $u_{t}$ with a zero memory parameter. Finally, Gaussianity is required for the sake of simplicity of the proofs as in Sun and Phillips (2003). Gaussianity of the noise precludes the possibility of LMSV where the noise is not Gaussian. However, Gaussianity of RV series has been supported by Andersen et al. (2003) and Lieberman and Phillips (2008) among others. LMSV has been allowed by Deo and Hurvich (2001) and Hurvich and Soulier (2002) for the original log periodogram
regression when the $u_{t}$ is a white noise (note that we allow for a weak dependent $u_{t}$ ) that is not accounted for in the estimation procedure. Gaussianity of $u_{t}$ is more difficult to relax if the noise is included in the estimation procedure since in that case a non linear transformation of its periodogram has to be considered for the asymptotics of the estimator. Gaussianity of the signal is even more difficult to avoid and only Velasco (2000) for the original log periodogram regression in a fully observable series has replaced that assumption by restrictions in higher order moments, but even in that simplest context the asymptotic properties are only obtained if tapering is previously applied, with the consequent loss of efficiency.

Assumption 3 imposes a local behaviour of the (pseudo) cross spectral density of signal and noise. We call $C_{y u}$ and $G_{y u}$ the low and high frequency correlation parameters respectively since the latter is multiplied by the frequency and is thus negligible with respect to the low frequency correlation for frequencies close to zero. It is based on the extended use of the fractional difference operator $(1-L)^{d}$ such that in its phasor form

$$
\begin{align*}
\left(1-e^{ \pm i \lambda}\right)^{d} & =\left(2 \sin \frac{\lambda}{2}\right)^{d} \exp \left\{ \pm i d\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)\right\}  \tag{3}\\
& =\lambda^{d}\left(1+O\left(\lambda^{2}\right)\right) \exp \left\{ \pm i d\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)\right\}
\end{align*}
$$

Hurvich et al. (2005) consider only the cosine term, ignoring the sine imaginary part so that only the low frequency correlation is accounted for. In some situations this omission generates a bias that would limit the number of frequencies or bandwidth used in their estimation further than claimed (see Remarks 3 and 4 below). We found it necessary also to include the sine imaginary part to fully account for the correlation of signal and noise.

Assumption 3 implies that the phase at zero frequency is fixed to be $d_{0} \pi / 2$, which is related to the use of the one sided fractional filter $(1-L)^{d}$, but other semiparametric estructures give also rise to such a phase (see Robinson 2008 and the examples below). We found this assumption necessary for identifiability of the low and high frequency correlation parameters. Consider instead an unknown phase parameter $\gamma_{0}$ and modify correspondingly Assumption 3 as in Robinson (2008) such that, as $\lambda \rightarrow 0^{+}$

$$
\begin{equation*}
\operatorname{Re}\left(f_{y u}(\lambda)\right)=\lambda^{-d_{0}}\left(C_{y u} \cos \gamma_{0}+G_{y u} \lambda \sin \gamma_{0}+O\left(\lambda^{\beta_{3}}\right)\right) . \tag{4}
\end{equation*}
$$

In this framework $C_{y u}$ and $\gamma_{0}$ are not jointly identifiable (the effect is similar to an unknown memory of the added noise discussed in Remark 10 below) and knowledge of one of them is
needed to allow for the estimation of the other since the information on both parameters is concentrated in a single term $C_{y u} \cos \gamma_{0}$. Assumption 3 imposes a particular behaviour of the phase not only at zero frequency but also at neighbouring frequencies. Another possibility is to specify the real part of the (pseudo) cross spectral density as in (4) with $d_{0} \pi / 2$ instead of $\gamma_{0}$ but in this case $G_{y u}$ loses its correlation interpretation because it includes a term depending on $d_{0}$ and generated by the approximation of the cosine.

Different spectral smoothness parameters are permitted in the (pseudo) spectral and cross spectral densities of $y_{t}$ and $u_{t}$. The restrictions $\beta_{1}>1+d_{0}, \beta_{2}>\beta_{1}-2 d_{0}$ and $\beta_{3}>\beta_{1}-d_{0}$ guarantee that the asymptotic bias of the proposed estimator is $O\left(\lambda^{\beta_{1}}\right)$ and also ensure identifiability of all the parameters to be estimated. For that $\beta_{1}>1+d_{0}$, $\beta_{2}>1-d_{0}$ and $\beta_{3}>1$ because otherwise the remainder in our regression model would be of an order of magnitude larger than $O\left(\lambda^{1+d_{0}}\right)$ and at least $G_{y u}$ would not be identifiable. These conditions are not very restrictive because in a wide range of situations, such as in ARMA and ARFIMA models, $\beta_{1}=\beta_{2}=\beta_{3}=2$.

For the purpose of illustration consider the two following cases (non stationary signals can be defined as explained in the Introduction).

Case 1: Let

$$
\begin{aligned}
y_{t} & =(1-L)^{-d_{0}} \alpha(L) w_{t} \quad \text { for } w_{t} \sim N I D\left(0, \sigma_{w}^{2}\right) \\
u_{t} & =\beta(L) \varepsilon_{t} \quad \text { for } \varepsilon_{t} \sim \operatorname{NID}\left(0, \sigma_{\varepsilon}^{2}\right) \\
E\left(w_{t} \varepsilon_{t}\right) & =\rho \quad \text { and } \quad E\left(w_{t} \varepsilon_{s}\right)=0 \quad t \neq s
\end{aligned}
$$

where $L$ is the lag operator, $\alpha(L)=\sum_{j=0}^{\infty} \alpha_{j} L^{j}$ and $\beta(L)=\sum_{j=0}^{\infty} \beta_{j} L^{j}$ such that $\alpha(x)$ is twice differentiable with the second derivative satisfying a Lipschitz condition of order $\iota$ at one (implying $\beta_{1}=2$ ) and $\beta(x)$ is twice differentiable with a bounded second derivative in the neighbourhood of one (such that $\beta_{2}=2$ ). For example $u_{t}$ can be a weak dependent ARMA and $y_{t}$ a stationary fractional ARIMA, but other models such as the Bloomfield exponential and fractionally integrated Bloomfield exponential are also covered. Here $f_{u}(0)=\sigma_{\varepsilon}^{2} \beta(1)^{2} / 2 \pi, C_{y}=\sigma_{w}^{2} \alpha(1)^{2} / 2 \pi, G_{y}=d_{0} / 12+\left\{\alpha^{\prime}(1)^{2}-\left[\alpha^{\prime \prime}(1)+\alpha^{\prime}(1)\right] \alpha(1)\right\} / \alpha(1)^{2}$ and the cross spectral density is

$$
f_{y u}(\lambda)=\frac{\rho}{2 \pi}\left(1-e^{-i \lambda}\right)^{-d_{0}} \alpha\left(e^{-i \lambda}\right) \overline{\beta\left(e^{-i \lambda}\right)}
$$

where the overline denotes complex conjugation. Using (3) we have that $f_{y u}(\lambda)$ satisfies Assumption 3 with $\beta_{3}=2$ and

$$
C_{y u}=\frac{\rho}{2 \pi} \alpha(1) \beta(1), G_{y u}=\frac{\rho}{2 \pi}\left[\alpha^{\prime}(1) \beta(1)-\beta^{\prime}(1) \alpha(1)\right]
$$

A particular situation arises when $y_{t}$ is a fractional noise and $u_{t}$ is white noise such that $\alpha(L)=\beta(L)=1$ and Assumptions 1 and 2 hold with $f_{u}(0)=\sigma_{\varepsilon}^{2} / 2 \pi, C_{y}=\sigma_{w}^{2} / 2 \pi$ and $G_{y}=d_{0} / 12$. Regarding the cross spectral density, Assumption 3 holds with $C_{y u}=\rho / 2 \pi$ and $G_{y u}=0$. More generally, $G_{y u}=0$ whenever $\alpha^{\prime}(1) \beta(1)-\beta^{\prime}(1) \alpha(1)=0$ and the imaginary part of $\left(1-e^{i \lambda}\right)^{-d}$ plays no role in the local specification of the cross spectral density such that the high frequency correlation in Assumption 3 can be discarded. This particular case represents the context covered by the extension of the local Whittle estimator in Hurvich et al. (2005).

Case 2: It is sometimes of interest to permit non contemporaneous correlation of signal and noise innovations such that

$$
E\left(w_{t} \varepsilon_{t-s}\right)=\rho \quad \text { and } \quad E\left(w_{t} \varepsilon_{t-k}\right)=0 \quad k \neq s
$$

for some $s \neq 0$, to allow for example for lagged adjustment to information. In this case the cross spectral density function is

$$
f_{y u}(\lambda)=\frac{\rho}{2 \pi} e^{-i \lambda s}\left(1-e^{-i \lambda}\right)^{-d_{0}} \alpha\left(e^{-i \lambda}\right) \overline{\beta\left(e^{-i \lambda}\right)}
$$

which satisfies

$$
\operatorname{Re}\left(f_{y u}(\lambda)\right)=\lambda^{-d_{0}}\left\{C_{y u}^{*} \cos \left[d_{0}\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)+\lambda s\right]+G_{y u}^{*} \lambda \sin \left[d_{0}\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)+\lambda s\right]+O\left(\lambda^{2}\right)\right\}
$$

where

$$
C_{y u}^{*}=\frac{\rho}{2 \pi} \alpha(1) \beta(1), \quad G_{y u}^{*}=\frac{\rho}{2 \pi}\left(\alpha^{\prime}(1) \beta(1)-\beta^{\prime}(1) \alpha(1)\right)
$$

Since

$$
\begin{aligned}
& \cos \left[d_{0}\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)+\lambda s\right]=\cos \left[d_{0}\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)\right]-\lambda s \sin \left[d_{0}\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)\right]+O\left(\lambda^{2}\right) \\
& \sin \left[d_{0}\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)+\lambda s\right]=\sin \left[d_{0}\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)\right]+\lambda s \cos \left[d_{0}\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)\right]+O\left(\lambda^{2}\right)
\end{aligned}
$$

$f_{y u}(\lambda)$ satisfies Assumption 3 with

$$
C_{y u}=C_{y u}^{*} \quad \text { and } \quad G_{y u}=G_{y u}^{*}-s C_{y u}
$$

The intertemporal correlation precludes the possibility of a null $G_{y u}$ even when $y_{t}$ is a fractional noise and $u_{t}$ is white noise.

## 3 ESTIMATION UNDER CORRELATION OF SIGNAL AND NOISE

Under Assumptions 1, 2 and 3, the (pseudo) spectral density of $z_{t}$ at frequency $\lambda_{j}$ satisfies

$$
\begin{equation*}
f_{z}\left(\lambda_{j}\right)=C_{y} \lambda_{j}^{-2 d_{0}}\left(1+\theta_{0}^{\prime} X_{j}\left(d_{0}\right)+\varsigma_{j}\right) \tag{5}
\end{equation*}
$$

where $\theta_{0}=\left(\theta_{10}, \theta_{20}, \theta_{30}\right)^{\prime}, X_{j}(d)=\left(A_{j}(d) \lambda_{j}^{d}, \lambda_{j}^{2 d}, B_{j}(d) \lambda_{j}^{1+d}\right)^{\prime}, A_{j}(d)=\cos \left[d\left(\frac{\lambda_{j}}{2}-\frac{\pi}{2}\right)\right]$, $B_{j}(d)=\sin \left[d\left(\frac{\lambda_{j}}{2}-\frac{\pi}{2}\right)\right], \varsigma_{j}=O\left(\lambda_{j}^{\beta_{1}}\right)$ under Assumption 1a) and $\varsigma_{j}=G_{y} \lambda_{j}^{\beta_{1}}+O\left(\lambda_{j}^{\beta_{1}+l}\right)$ for $l=\min \left(\iota, d_{0}\right)$ under Assumption 1b) and

$$
\theta_{10}=\frac{2 C_{y u}}{C_{y}}, \quad \theta_{20}=\frac{f_{u}(0)}{C_{y}}, \quad \theta_{30}=\frac{2 G_{y u}}{C_{y}}
$$

Taking logarithms of (5) and considering only Fourier frequencies $\lambda_{j}=2 \pi j / n, j=$ $1,2, \ldots[n / 2]$ for $n$ the sample size, we have

$$
\begin{equation*}
\log I_{z j}=a_{0}+d_{0}\left(-2 \log \lambda_{j}\right)+\log \left(1+\theta_{0}^{\prime} X_{j}\left(d_{0}\right)\right)+\varsigma_{j}+U_{z j} \tag{6}
\end{equation*}
$$

where $a_{0}=\log C_{y}-c, c=0,577216$ is Euler's constant, $U_{z j}=\log \left(I_{z j} f_{z}^{-1}\left(\lambda_{j}\right)\right)+c$ and $I_{z j}$ is the periodogram of $z_{t}, t=1,2, \ldots, n$, at frequency $\lambda_{j}$

$$
I_{z j}=I_{z}\left(\lambda_{j}\right)=\left|w_{z j}\right|^{2} \quad \text { for } \quad w_{z j}=w_{z}\left(\lambda_{j}\right)=\frac{1}{\sqrt{2 \pi n}} \sum_{t=1}^{n} z_{t} \exp \left(-i \lambda_{j} t\right)
$$

The $\varsigma_{j}$ in (6) is different from that in (5) but we use the same notation because they are asymptotically equivalent in the sense that they coincide up to an $o\left(\lambda_{j}^{\beta_{1}}\right)$ term.

Our main interest is to estimate the memory parameter $d_{0}$ of the latent signal $y_{t}$, although the rest of parameters $\theta_{10}, \theta_{20}$ and $\theta_{30}$ may also play an important role and their joint estimation not only reduces the bias of the estimates of $d_{0}$ but is also of interest in itself. The parameter $\theta_{20}$ is the long run noise to signal ratio and $\theta_{10}$ and $\theta_{30}$ correspond to the low frequency and high frequency correlations respectively between signal and noise. Any correlation between innovations of signal and noise entails $\theta_{10} \neq 0$. It also generally implies $\theta_{30} \neq 0$ but there are some particular cases where $\theta_{30}=0$, as indicated in the previous section. Non contemporaneous correlation, however, always implies $\theta_{30} \neq 0$.

In order to avoid possible inconsistencies caused by spectral misspecifications at frequencies far from the origin, we focus only on the $m$ Fourier frequencies closest to zero, as is usual
in other semiparametric or local estimators of $d_{0}$. The Augmented Log-Periodogram regression Estimator (ALPE) is obtained by applying least squares to the non linear regression model

$$
\begin{equation*}
\log I_{z j}=a+d\left(-2 \log \lambda_{j}\right)+\log \left(1+\theta^{\prime} X_{j}(d)\right)+U_{z j} \quad j=1,2, \ldots, m, \tag{7}
\end{equation*}
$$

where the $\varsigma_{j}$ term that is omitted will lead the bias of the estimates. This is an extension of the ALPE of Arteche (2006) to account for possible correlation between $y_{t}$ and $u_{t}$ covering also non stationary values of $d_{0}$. Approximating $\log \left(1+\theta^{\prime} X_{j}(d)\right)$ locally by $\theta^{\prime} X_{j}(d)$ as in Sun and Phillips (2003) results in a regression model that is linear on $\theta$ but still non linear on $d$, giving rise to an extra bias term of order $O\left(\lambda_{m}^{2 d_{0}}\right)$ (Arteche, 2006) and making $\theta_{20}$ and $\theta_{30}$ unidentifiable if signal and noise are correlated. Also, ignoring the high frequency correlation as in Hurvich et al. (2005) would create an extra bias term of order $O\left(\lambda_{m}^{1+d_{0}}\right)$ except in those particular cases where $\theta_{30}=0$ mentioned in the previous section- that would limit the size of $m$, requiring a more restrictive bandwidth than that allowed in Assumption 4 below. In consequence the assumption in equation (3.9) in Hurvich et al. (2005) seems insufficient to make the effect of this omitted term asymptotically negligible and a more restrictive assumption on the evolution of $m$ seems necessary as mentioned in Remark 1 below and acknowledged by the authors in their corrigendum (Hurvich et al., 2008).

The ALPE is formally defined as

$$
\begin{equation*}
\left(\hat{d}_{A L P}, \hat{\theta}_{A L P}\right)=\arg \min _{\Delta \times \Theta} Q(d, \theta) \tag{8}
\end{equation*}
$$

where $\Delta=\left[\Delta_{1}, \Delta_{2}\right], 0<\Delta_{1}<\Delta_{2}<1, \Theta=\Theta_{1} \times \Theta_{2} \times \Theta_{3}$ for $\Theta_{1}=\left[\Theta_{11}, \Theta_{12}\right],-\infty<$ $\Theta_{11}<\Theta_{12}<\infty, \Theta_{2}=\left[0, \Theta_{22}\right], 0<\Theta_{22}<\infty, \Theta_{3}=\left[\Theta_{31}, \Theta_{32}\right],-\infty<\Theta_{31}<\Theta_{32}<\infty$, and

$$
Q(d, \theta)=\frac{1}{m} \sum_{j=1}^{m}\left\{\log ^{\dagger} I_{z j}+d\left(2 \log \lambda_{j}\right)^{\dagger}-\log ^{\dagger}\left(1+\theta^{\prime} X_{j}(d)\right)\right\}^{2}
$$

where for a general $\xi_{j}$ we use the notation $\xi_{j}^{\dagger}=\xi_{j}-\bar{\xi}$ where $\bar{\xi}=\sum \xi_{j} / m$.

## 4 ASYMPTOTIC PROPERTIES OF THE ALPE

Theorem 1 Under Assumptions 1a)-3, $\hat{d}_{A L P}-d_{0}=o_{p}(1)$ if $m^{-1}+m n^{-1} \rightarrow 0$ as $n \rightarrow$ $\infty$ and $\hat{d}_{A L P}-d_{0}=O_{p}\left(\lambda_{m}^{1+d_{0}}\right), \hat{\theta}_{1 A L P}-\theta_{10}=O_{p}\left(\lambda_{m}\right), \hat{\theta}_{2 A L P}-\theta_{20}=o_{p}\left(\lambda_{m}^{1-d_{0}}\right)$ and $\hat{\theta}_{3 A L P}-\theta_{30}=o_{p}(1)$ if $m n^{-1}+n^{2\left(1+d_{0}\right)(1+\delta)} m^{-2\left(1+d_{0}\right)(1+\delta)-1} \rightarrow 0$ as $n \rightarrow \infty$ for some arbitrarily small $\delta>0$ and $d_{0}<3 / 4$.

Theorem 1 shows the consistency of $\hat{d}_{A L P}$ as long as $0<d_{0}<1$. For the consistency of $\hat{\theta}_{A L P}$ we need a more refined rate of convergence of $\hat{d}_{A L P}$ to avoid the flatness of $Q(d, \theta)$ as a function of $\theta$. This is only achieved for $d_{0}<3 / 4$, when the $m^{2 d_{0}-2}$ terms in the bounds in Corollary 2 are dominated by the $m^{-1 / 2}$ terms. With that bounds we get consistency (at different rates) of the estimators of all the parameters.

For the asymptotic normality of $\left(\hat{d}_{A L P}, \hat{\theta}_{A L P}\right)$ a more restrictive assumption on the rate of increase of the bandwidth is required.

Assumption 4(K): For $\delta>0$ arbitrarily small and $0<d_{0}<3 / 4$, as $n \rightarrow \infty$,

$$
\frac{n^{2\left(d_{0}+1\right)(1+\delta)}}{m^{1+2\left(d_{0}+1\right)(1+\delta)}} \rightarrow 0 \quad \text { and } \quad \frac{m^{\beta_{1}+1 / 2}}{n^{\beta_{1}}} \rightarrow K
$$

for a finite constant $K$.

The first condition in Assumption 4(K) imposes a lower bound on the growth rate of $m$, ensuring the consistency of the ALPE. A larger bandwidth is required as the value of $d_{0}$ increases to guarantee consistency of all the estimators of the elements in $\theta$. It ensures that all the elements in the diagonal of the normalizing matrix $D_{n}$ of the gradient and Hessian defined in Theorem 2 go to infinity. The upper bound is imposed by the second condition and is the conventional $O\left(n^{2 \beta_{1} /\left(2 \beta_{1}+1\right)}\right)$ rate in the log periodogram regression when applied to a fully observable long memory series. These two restrictions are always compatible because $\beta_{1}>d_{0}+1$ and $\delta$ is arbitrarily small.

The asymptotic distribution depends on the location of $\theta_{0}$ in the parameter space. We first consider the case of existence of added noise such that all the parameters to be estimated are in the interior of the parameter set.

Assumption 5: $\left(d_{0}, \theta_{0}\right)=\left(d_{0}, \theta_{10}, \theta_{20}, \theta_{30}\right)$ is an interior point of the parameter space, $\left(d_{0}, \theta_{10}, \theta_{20}, \theta_{30}\right) \in\left(\Delta_{1}, \Delta_{2}\right) \times\left(0, \Theta_{11}\right) \times\left(\Theta_{21}, \Theta_{22}\right) \times\left(\Theta_{31}, \Theta_{32}\right)$.

Theorem 2 Under Assumptions 1b)-4(K) and 5

$$
D_{n}\binom{\hat{d}_{A L P E}-d_{0}}{\hat{\theta}_{A L P}-\theta_{0}} \xrightarrow{d} N\left(\Omega^{-1} b, \frac{\pi^{2}}{6} \Omega^{-1}\right)
$$

for $D_{n}=D_{n}\left(d_{0}\right)$ and $D_{n}(d)=\sqrt{m} \operatorname{diag}\left(1, \lambda_{m}^{d} \cos (\pi d / 2), \lambda_{m}^{2 d}, \lambda_{m}^{1+d} \sin (\pi d / 2)\right), b=b\left(d_{0}\right)$
with

$$
b(d)=\left(\begin{array}{c}
-\frac{2 \beta_{1}}{\left(\beta_{1}+1\right)^{2}} \\
\frac{\begin{array}{c}
\beta_{1} d
\end{array}}{\frac{\left(\beta_{1}+1\right)\left(\beta_{1}+1\right)(d+1)}{2 \alpha_{1}}} \\
\frac{2 d_{1}}{\left(\beta_{1}+1\right)\left(2 d+\beta_{1}+1\right)(2 d+1)} \\
-\frac{\beta_{1}(1+d)}{\left(\beta_{1}+1\right)\left(d+\beta_{1}+2\right)(d+2)}
\end{array}\right) K(2 \pi)^{\beta_{1}} G_{y}
$$

and $\Omega=\Omega\left(d_{0}\right)$ for

$$
\Omega(d)=\left(\begin{array}{cccc}
4 & -\frac{2 d}{(1+d)^{2}} & -\frac{4 d}{(1+2 d)^{2}} & \frac{2(1+d)}{(2+d)^{2}} \\
& \frac{d^{2}}{(2 d+1)(1+d)^{2}} & \frac{2 d^{2}}{(2 d+1)(1+d)(3 d+1)} & -\frac{4 d^{2}}{2(d+1)(2+d)} \\
& & \frac{2 d(d+1)}{(4 d+1)(1+2 d)^{2}} & -\frac{2 d(2 d)}{(3 d+2)(2 d+1)(d+2)} \\
& & & \frac{1+d)^{2}}{(2 d+3)(2+d)^{2}}
\end{array}\right)
$$

Remark 1: Theorems 1 and 2 show that the ALPE is consistent for $0<d_{0}<1$ and asymptotically normal as long as $d_{0}<3 / 4$. Lemmas 3 and 6 in Appendix B can similarly be used to show consistency (for $d_{0}<1$ ) and asymptotic normality (for $d_{0}<3 / 4$ ) of the estimators proposed by Sun and Phillips (2003) and Arteche (2006) when signal and noise are independent, extending the work by Velasco (1999) for nonstationary series. Note that Velasco required to trim out $l$ frequencies close to the origin. However, bounding the contributions of the low frequencies as advocated by Hurvich et al. (1998) and used in the proof of Theorem 2 results in that this trimming is not necessary for the classical log periodogram regression nor the extensions here considered.

Remark 2: The inclusion of regressors in the estimation procedure inflates the asymptotic variance of the estimator of $d_{0}$ by a multiplicative constant that is higher the lower $d_{0}$ is. Table 1 shows the asymptotic variances of $\sqrt{m} \hat{d}_{A L P}, \sqrt{m} \hat{d}_{1}$ and $\sqrt{m} \hat{d}_{2}$ for $d_{0}=0.2,0.36$ and 0.48 , where $\hat{d}_{1}$ is obtained by nonlinear log-periodogram regression with no correlation between signal and noise (that is omitting $\theta_{1}$ and $\theta_{3}$ in the estimation procedure, see Sun and Phillips, 2003 and Arteche, 2006) and $\hat{d}_{2}$ ignores the high frequency correlation (omits $\theta_{3}$ ) as in Hurvich et al. (2005). The inclusion of these regressors reduces the bias and allows a broader bandwidth such that, although the variance increases significantly, the asymptotic efficiency can be improved by using a larger $m$.

Remark 3: If $\theta_{30}=0$ the asymptotic covariance matrix of the correspondingly restricted ALPE is $\pi^{2} / 6$ times the inverse of the $(3 \times 3)$ left upper submatrix of $\Omega$. This is the covariance matrix of the local Whittle estimator of Hurvich et al. (2005), denoted by $\Gamma^{*}$ in their Proposition 4.1 ${ }^{1}$, apart from the multiplicative constant, which is $\pi^{2} / 6$ here

[^1]Table 1: Asymptotic variances of $\sqrt{m} \hat{d}_{1}, \sqrt{m} \hat{d}_{2}$ and $\sqrt{m} \hat{d}_{A L P}$

|  | $\hat{d}_{1}$ | $\hat{d}_{2}$ | $\hat{d}_{A L P}$ |
| :---: | :---: | :---: | :---: |
| $d_{0}=0.2$ | 5.038 | 181.354 | 609.551 |
| $d_{0}=0.36$ | 2.347 | 33.493 | 100.855 |
| $d_{0}=0.48$ | 1.714 | 16.297 | 45.759 |

due to the different estimation procedure. This is the typical discrepancy between other local Whittle and log-periodogram regression based estimators. However, if $\theta_{30} \neq 0$, the omission of $\theta_{3}$ in the restricted estimation generates an extra bias due to the $O\left(\lambda_{j}^{1+d_{0}}\right)$ high frequency correlation component in the (pseudo) cross spectral density, which is not explicitly considered in the estimation. A more restrictive assumption on the growth rate of $m$ should then be imposed for asymptotic normality. Precisely $m^{2\left(1+d_{0}\right)+1} n^{-2\left(d_{0}+1\right)}=O(1)$ as $n \rightarrow \infty$ should hold instead of the second part of Assumption 4. A similar condition seems also to be necessary in Hurvich et al. (2005) (see the corrigendum by Hurvich et al. 2008).

Remark 4: The analogue local Whittle type estimator of our ALPE is the Modified Gaussian Semiparametric Estimator (MGSE) defined as

$$
\begin{equation*}
(\tilde{d}, \tilde{\theta})=\arg \min _{\Delta \times \Theta}\left\{\log \left(\frac{1}{m} \sum_{j=1}^{m} \frac{\lambda_{j}^{2 d} I_{z j}}{1+\theta^{\prime} X_{j}(d)}\right)+\frac{1}{m} \sum_{j=1}^{m} \log \left\{\lambda_{j}^{-2 d}\left(1+\theta^{\prime} X_{j}(d)\right)\right\}\right\} \tag{9}
\end{equation*}
$$

which corresponds to the estimator of Hurvich et al. (2005) but including the high frequency correlation. This case represents a new parametrization (P3) in Hurvich et al. (2005) and under a similar set of assumptions their results cover also this possibility such that consistency and asymptotic normality are expected to hold with an asymptotic covariance matrix as that in Theorem 2 with $\pi^{2} / 6$ replaced by one. Comparing their assumptions with those needed here for the ALPE, the main advantage of the local Whittle extension is that Gaussianity of signal and noise are not required and instead it is assumed that the stationary part of the signal admits an infinite moving order representation with respect to a martingale difference sequence with bounded fourth moment and $u_{t}$ is a zero mean white noise with finite fourth moment. The assumptions required for the short memory part of the signal are implied by our Assumption 1. Note also that the ALPE covers a more general situation because it allows for weak dependence in the added noise and by two and $\cos \left(d_{0} \pi / 2\right)$ should be multiplying in their normalizing matrix as in our $D_{n}$.
non contemporaneous correlation. Finally the upper bound in the rate of increase of the bandwidth for the asymptotic normality of the MGSE is slightly more restrictive due to a $\log m$ term that appears in formula (3.9) of Hurvich et al (2005) and is avoided in our Assumption 4.

From an empirical perspective, the main advantage of the standard LPE over the local Whittle lies in its simple implementation as a linear regression. This property however is lost in the ALPE as well as in other log periodogram regression based estimators that account for the added noise (Sun and Phillips, 2003, Arteche, 2006), which require nonlinear optimization. Nevertheless the ALPE seems to be superior in finite samples to the MGSE, at least in the cases analyzed in the Monte Carlo in Section 7. This can be partly explained by the fact that the empirical implementation of both estimators need to restrict $1+\theta^{\prime} X_{j}(d)$ to be positive and this expression appears twice in the contrast function of the MGSE and only once in the ALPE.

Remark 5: If the (pseudo) cross spectral density is not considered explicitly in the estimation procedure when in fact $\theta_{10} \neq 0$, the wrongly restricted estimator of $d_{0}$ remains consistent with an appropriate bandwidth choice but $\theta_{20}$ is not identifiable because the remainder is of order $O\left(\lambda_{m}^{d_{0}}\right)$, i.e. higher than the order of the regressor corresponding to $\theta_{20}$. In consequence $\theta_{20}$ can not be estimated consistently. Based on this characteristic, a Hausman type test for correlation between signal and noise is introduced in Section 5.

Remark 6: If the high frequency correlation is not included in the estimation when actually $\theta_{30} \neq 0$, the restricted estimators of $d_{0}, \theta_{10}$ and $\theta_{20}$ are consistent if $n^{4 d_{0}(1+\delta)} m^{-1-4 d_{0}(1+\delta)} \rightarrow$ 0 for some small $\delta>0$. The asymptotic bias of the restricted ALPE of $d_{0}$ can be approximated in this case by $\theta_{30} \lambda_{m}^{1+d_{0}} \sin \left(\pi d_{0} / 2\right) H\left(d_{0}\right)$ where $H\left(d_{0}\right)>0$ is a constant function of $d_{0}$. The same expression approximates the asymptotic bias of the local Whittle estimator of Hurvich et al. (2005).

Remark 7: The asymptotic bias of $\hat{d}_{A L P}$ can be approximated by

$$
\operatorname{Abias}\left(\hat{d}_{A L P}\right)=\left(\frac{m}{n}\right)^{\beta_{1}} \tilde{\Omega}_{1} b_{K}
$$

where $\tilde{\Omega}_{1}$ is the first row of $\Omega^{-1}$ and $b_{K}=b / K$. The asymptotic variance is

$$
\operatorname{Avar}\left(\hat{d}_{A L P}\right)=\frac{\pi^{2}}{6 m} \tilde{\Omega}_{11}
$$

and consequently the asymptotic mean squared error can be approximated by

$$
\operatorname{AMSE}\left(\hat{d}_{A L P}\right)=\frac{\pi^{2}}{6 m} \tilde{\Omega}_{11}+\left(\frac{m}{n}\right)^{2 \beta_{1}}\left(\tilde{\Omega}_{1} b_{K}\right)^{2}
$$

The "optimal" bandwidth that minimizes $A M S E\left(\hat{d}_{A L P}\right)$ is

$$
m_{A L P E}^{o p t}=n^{2 \beta_{1} /\left(2 \beta_{1}+1\right)}\left[\frac{\pi^{2} \tilde{\Omega}_{11}}{24\left(\tilde{\Omega}_{1} b_{K}\right)^{2}}\right]^{1 /\left(2 \beta_{1}+1\right)}
$$

such that $\operatorname{AMSE}\left(\hat{d}_{A L P}\right)=O\left(n^{-2 \beta_{1} /\left(2 \beta_{1}+1\right)}\right)$ if $m=m_{A L P E}^{o p t}$. The ALPE then achieves the same rate of mean square error convergence as the standard LPE applied to a fully observable long memory series.

Remark 8: The bandwidths that minimize the AMSE of $\hat{\theta}_{i, A L P E}$ are $O\left(n^{2 \beta_{1} /\left(2 \beta_{1}+1\right)}\right)$ for all $i=1,2,3$, although with different multiplicative constants in each case. The optimal rates of convergence of $\hat{\theta}_{i, A L P E}$ in a mean squared error sense are then $O\left(n^{-\left(\beta_{1}-d_{0}\right) /\left(2 \beta_{1}+1\right)}\right)$, $O\left(n^{-\left(\beta_{1}-2 d_{0}\right) /\left(2 \beta_{1}+1\right)}\right)$ and $O\left(n^{-\left(\beta_{1}-1-d_{0}\right) /\left(2 \beta_{1}+1\right)}\right)$ for $i=1,2,3$ respectively, which can be quite slow for large $d_{0}$.

Remark 9: We do not consider the possibility of $d_{0}=0$, so that the asymptotic distribution in Theorem 2 can not be used to test the hypothesis of short memory. However if $d_{0}=0$ there is not need to extend the original LPE and GSE since such an extension does not imply a bias reduction. The classical LPE and GSE can then be used for that purpose as suggested by Hurvich and Soulier (2002) and Hurvich et al. (2005) since both maintain the asymptotic properties as if no added noise were present. In particular both are asymptotically normal and $n^{\beta_{1} /\left(2 \beta_{1}+1\right)}$-consistent with an optimal bandwidth choice. Note also that if $d_{0}=0$ the vector of parameters $\theta$ need not be identifiable. Thus, if we suspect that the value of $d_{0}$ could be zero, we should test that possibility by means of the standard LPE or GSE. If we find evidence of a positive $d_{0}$ then some extension accounting for the possible existence of noise should be applied. The results in Theorem 2 can be used however to test the stationarity of the series by means of the null hypothesis $d_{0}=1 / 2$ against $d_{0}<1 / 2$.

Remark 10: The literature on perturbed long memory has focused on a weak dependent added noise, which is the main case of interest in economics and finance, as illustrated in the introduction. Consider however that, instead of Assumption 2, $u_{t}$ is a long memory process with a spectral density satisfying

$$
f_{u}(\lambda)=C_{u} \lambda^{-2 d_{u 0}}\left(1+O\left(\lambda^{\beta_{2}}\right)\right)
$$

as $\lambda \rightarrow 0^{+}$, for a positive finite constant $C_{u}$ and $-0.5<d_{u 0}<d_{0}$ where $d_{u 0}$ is a new parameter to be estimated. In this context the (pseudo) cross spectral density function of $y_{t}$ and $u_{t}$ satisfies as $\lambda \rightarrow 0^{+}$,
$\operatorname{Re}\left(f_{y u}(\lambda)\right)=\lambda^{-d_{0}-d_{u 0}}\left(C_{y u} \cos \left[\left(d_{0}-d_{u 0}\right)\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)\right]+G_{y u} \lambda \sin \left[\left(d_{0}-d_{u 0}\right)\left(\frac{\lambda}{2}-\frac{\pi}{2}\right)\right]+O\left(\lambda^{\beta_{3}}\right)\right)$
for finite constants $C_{y u}$ and $G_{y u}$. The two examples given in Section 2 with the added noise replaced by $u_{t}=(1-L)^{-d_{u 0}} \beta(L) \varepsilon_{t}$ satisfy this specification. The (pseudo) spectral density function of $z_{t}$ is then

$$
f_{z}\left(\lambda_{j}\right)=C_{y} \lambda_{j}^{-2 d_{0}}\left(1+\theta_{0}^{\prime} X_{j}\left(d_{0}-d_{u 0}\right)+\varsigma_{j}\right)
$$

and the regression model is

$$
\begin{equation*}
\log I_{z j}=a+d\left(-2 \log \lambda_{j}\right)+\log \left(1+\theta^{\prime} X_{j}\left(d-d_{u}\right)\right)+U_{z j} \quad j=1,2, \ldots, m \tag{10}
\end{equation*}
$$

There is here a problem of asymptotic identification. Denote by $\psi_{0}=\left(d_{0}, \theta_{0}^{\prime}, d_{u 0}\right)^{\prime}$ the parameters to be estimated and let

$$
q_{j}(\psi)=d\left(-2 \log \lambda_{j}\right)^{\dagger}+\log \left(1+\theta^{\prime} X_{j}\left(d-d_{u}\right)\right)^{\dagger}
$$

Denote by $X_{0}$ the $m \times 5$ matrix with elements $\left[X_{0}\right]_{j i}=\partial q_{j}\left(\psi_{0}\right) / \partial \psi_{i}$ and $D_{n}^{\perp}=D_{n}^{\perp}\left(d_{0}-d_{u 0}\right)$ where $D_{n}^{\perp}(d)=\sqrt{m} \operatorname{diag}\left(1, \lambda_{m}^{d} \cos (\pi d / 2), \lambda_{m}^{2 d}, \lambda_{m}^{1+d} \sin (\pi d / 2), \lambda_{m}^{d} \cos (\pi d / 2) \log \lambda_{m}\right)$. Approximating sums by integrals it can be shown that

$$
\lim _{n \rightarrow \infty} D_{n}^{\perp-1} X_{0}^{\prime} X_{0} D_{n}^{\perp-1}=\Omega^{\perp}
$$

where $\Omega^{\perp}=\Omega^{\perp}\left(d_{0}-d_{u 0}\right)$ such that $\Omega^{\perp}(d)$ is a $5 \times 5$ symmetric matrix with the first submatrix of order 4 equal to $\Omega(d)$ in Theorem 2 and the fifth column (row) equal to the vector

$$
\theta_{10}\left(\begin{array}{c}
\frac{2 d}{(1+d)^{2}} \\
-\frac{d^{2}}{(2 d+1)(1+d)^{2}} \\
-\frac{2 d^{2}}{(2 d+1)(1+d)(3 d+1)} \\
\frac{d}{2(d+1)(2+d)} \\
\frac{d^{2} \theta_{10}}{(2 d+1)(1+d)^{2}}
\end{array}\right)
$$

The matrix $\Omega^{\perp}$ is also the limit of the properly normalized Hessian matrix and is singular since the last column is just $-\theta_{10}$ times the second one, so the strong asymptotic identifiability condition is not satisfied (Davidson and MacKinnon, 2004, Chapter 6). Note that the
information of the regression model on $d_{u}$ is asymptotically dominated by $\theta_{1} A_{j}\left(d-d_{u}\right) \lambda_{j}^{d-d_{u}}$ such that both $d_{u}$ and $\theta_{1}$ cannot be identified.

A long memory added noise precludes identifiability in the estimation techniques that account for it, whereas it only affects the rate of convergence of the basic local Whittle and LPE estimators that ignore the noise by constraining the rate of increase of the bandwidth (Arteche, 2004). For identifiability $d_{u 0}$ should be known but this is not realistic unless $d_{u 0}=0$ as assumed before, where such an imposition relies on the characteristics of the noise in the different situations mentioned in the introduction.

In practice we are unlikely to be able to discern a priori whether the series is perturbed by an added noise or not. It is then interesting to analyze also the case of no added noise such that $\theta_{20}=0$ lies on the boundary of the parameter space, which affects the limiting distribution of the estimators. Note also that $\theta_{20}=0$ precludes the possibility of correlation since it obviously implies $\theta_{10}=\theta_{30}=0$.

Theorem 3 Let Assumption $1 b)-4(K)$ hold and let $d_{0} \in\left(\Delta_{1}, \Delta_{2}\right)$ and $\theta_{0}=0$. Then

$$
\begin{aligned}
& \sqrt{m}\left(\hat{d}_{A L P}-d_{0}\right) \quad \xrightarrow{d}-\tilde{\Omega}_{1} \eta\left\{\tilde{\Omega}_{3} \eta \leq 0\right\}-\Omega_{11}^{-1} \eta_{1}\left\{\tilde{\Omega}_{3} \eta>0\right\} \\
& D_{n}^{* *}\left(\hat{\theta}_{A L P}-\theta_{0}\right) \quad \xrightarrow{d}-\tilde{\Omega}^{* *} \eta\left\{\tilde{\Omega}_{3} \eta \leq 0\right\}
\end{aligned}
$$

where $\tilde{\Omega}_{i}$ is the $i$-th row of the matrix $\tilde{\Omega}=\Omega^{-1}, D_{n}^{* *}$ is the $3 \times 3$ low right submatrix of $D_{n}$, $\tilde{\Omega}^{* *}$ is the low $3 \times 4$ submatrix of $\Omega^{-1}$ and $\eta=\left(\eta_{1}, \eta_{2}, \eta_{3}, \eta_{4}\right)^{\prime} \sim N\left(-b, \pi^{2} \Omega / 6\right)$.

The proof of Theorem 3 is a straightforward extension of that in Theorem 4 in Sun and Phillips (2003) and is thus omitted.

## 5 TESTS FOR CORRELATION BETWEEN SIGNAL AND NOISE

The asymptotic distribution of the ALPE makes easy implementation of standard asymptotic inference conceivable not only on $d_{0}$ but also on the components of the vector $\theta_{0}$. In particular, it is of special interest to test the hypothesis of no correlation between signal and noise for two reasons. First, it is of interest per se because the existence of correlation
influences subsequent analysis, for example for forecasting RVs or for estimating economic mechanisms involving series with a measurement error correlated with the latent variable. Second, it is of technical interest because it can be used as a tool for prior selection of a suitable estimation strategy. If no evidence of correlation is found the ALPE and MGSE should be adapted to this information because in that case introducing terms to account for the correlation in the estimation procedure inflates unnecessarily the variance with the consequent loss of efficiency.

However, testing such a hypothesis is not a trivial issue since it involves assessing the correlation between two unobservable series whose spectral behaviour is only locally restricted. Uncorrelation between signal and noise corresponds in our local setup to the null hypothesis

$$
\begin{equation*}
H_{0}: \theta_{10}=\theta_{30}=0 \tag{11}
\end{equation*}
$$

The Wald test statistic for this hypothesis is

$$
W=\frac{6}{\pi^{2}} \hat{\theta}_{A L P}^{*} D_{n}^{*}\left(\hat{d}_{A L P}\right) \Omega^{*}\left(\hat{d}_{A L P}\right) D_{n}^{*}\left(\hat{d}_{A L P}\right) \hat{\theta}_{A L P}^{*}
$$

where $\hat{\theta}_{A L P}^{*}=\left(\hat{\theta}_{1, A L P}, \hat{\theta}_{3, A L P}\right)^{\prime}, D_{n}^{*}(d)=\sqrt{m}\left(\lambda_{m}^{d} \cos (\pi d / 2), \lambda_{m}^{1+d} \sin (\pi d / 2)\right)$ and $\Omega^{*}(d)=$ $\Omega_{1}(d)-\Omega_{3}(d) \Omega_{2}^{-1}(d) \Omega_{3}(d)^{\prime}$ where $\Omega_{1}, \Omega_{2}$ and $\Omega_{3}$ are $2 \times 2$ matrices with elements (by row) the $(2,2),(2,4),(4,2)$ and $(4,4)$ (for $\left.\Omega_{1}\right),(2,1),(2,3),(4,1)$ and $(4,3)$ (for $\left.\Omega_{3}\right)$ and $(1,1)$, $(1,3),(3,1)$ and $(3,3)$ (for $\Omega_{2}$ ) of the matrix $\Omega(d)$ in Theorem 2 . The following corollary establishes the typical properties of the Wald type testing procedure based on $W$. The proof is straightforward and is omitted.

Corollary 1 Let Assumptions 1b)-3, $\mathcal{4}(0)$ and 5 hold. Under $H_{0}$ in (11), $W \xrightarrow{d} \chi_{2}^{2}$. By contrast $W \xrightarrow{p} \infty$ if $\theta_{10} \neq$ and/or $\theta_{30} \neq 0$. Also under the local alternative $H_{1}:\left(\theta_{10}, \theta_{30}\right)^{\prime}=$ $D_{n}^{*-1}\left(d_{0}\right) \delta$ for a non null vector $\delta=\left(\delta_{1}, \delta_{3}\right)^{\prime}$, W has a non central chi-squared asymptotic distribution, $\chi_{2}^{2}\left(6 \delta^{\prime}\left(\pi^{2} \tilde{\Omega}^{*}\left(d_{0}\right)\right)^{-1} \delta\right)$.

Note that in the examples described in Section $2 \theta_{30} \neq 0$ implies $\theta_{10} \neq 0$ but $\theta_{30}=0$ and $\theta_{10} \neq 0$ is also a possibility under the alternative. The slow convergence of the estimators of $\theta_{10}$ and $\theta_{30}$ affects the finite sample performance of this test. In fact, we have found through simulations (not reported but available upon request) that the finite sample performance is rather poor. In order to avoid estimation of $\theta_{10}$ and $\theta_{30}$ an asymptotically equivalent

Lagrange Multiplier (LM) type test can be used. For the null hypothesis in (11) the LM test statistic has the form

$$
\begin{equation*}
L M=S^{*}\left(\hat{d}_{A L P}^{R}, \hat{\theta}_{A L P}^{R}\right)^{\prime} D_{n}^{*}\left(\hat{d}_{A L P}^{R}\right)^{-1}\left[\frac{\pi^{2}}{6} \Omega^{*}\left(\hat{d}_{A L P}^{R}\right)\right]^{-1} D_{n}^{*}\left(\hat{d}_{A L P}^{R}\right)^{-1} S^{*}\left(\hat{d}_{A L P}^{R}, \hat{\theta}_{A L P}^{R}\right) \tag{12}
\end{equation*}
$$

where $S^{*}(d, \theta)=\left(S_{2}(d, \theta), S_{4}(d, \theta)\right)$ is the vector of the second and fourth elements of the score $S(d, \theta)$ in the proof of Theorem 2 in Appendix A and $\hat{d}_{A L P}^{R}, \hat{\theta}_{A L P}^{R}$ are the ALPE under the restriction of the null. These estimators correspond to those analyzed in Arteche (2006) under independence of signal and noise. The asymptotic properties of the testing procedure based on this statistic are those of the Wald type test in Corollary 1. The matrix $\Omega^{*}\left(\hat{d}_{A L P}^{R}\right)$ can be replaced by consistent estimates. In particular we can form similar matrices with the elements of the Hessian estimates of $\Omega$ in the proof of Theorem 2, $D_{n}\left(\hat{d}_{A L P}^{R}\right)^{-1} H\left(\hat{d}_{A L P}^{R}, \hat{\theta}_{A L P}^{R}\right) D_{n}\left(\hat{d}_{A L P}^{R}\right)^{-1}$ or $D_{n}\left(\hat{d}_{A L P}^{R}\right)^{-1} J\left(\hat{d}_{A L P}^{R}, \hat{\theta}_{A L P}^{R}\right) D_{n}\left(\hat{d}_{A L P}^{R}\right)^{-1}$. Using these alternatives the LM statistic has a simpler expression since there is no need to use the normalizing matrix $D_{n}$ whose elements cancel out in (12), so no information is needed on the non standard rates of convergence of the estimators of the different parameters. The form of the statistic in this case is

$$
\begin{equation*}
L M=S^{*}\left(\hat{d}_{A L P}^{R}, \hat{\theta}_{A L P}^{R}\right)^{\prime}\left[\frac{\pi^{2}}{6} \Xi^{*}\left(\hat{d}_{A L P}^{R}, \hat{\theta}_{A L P}^{R}\right)\right]^{-1} S^{*}\left(\hat{d}_{A L P}^{R}, \hat{\theta}_{A L P}^{R}\right) \tag{13}
\end{equation*}
$$

with $\Xi^{*}$ defined similarly to $\Omega^{*}$ with respect to $\Xi=H$ or $J$.

Considering that under correlation between signal and noise $\theta_{20}$ cannot be identified unless the correlation is explicitly considered in the estimation procedure, a Hausman type test for correlation can be easily designed based on the difference between $\hat{\theta}_{2, A L P E}$, consistent under null and alternative hypotheses but less efficient than $\hat{\theta}_{2, A L P E}^{R}$ if no correlation exists, and $\hat{\theta}_{2, A L P E}^{R}$ which is not consistent under the alternative of correlated signal and noise.

Theorem 4 Under Assumptions 1b)-3, 4(0), 5 and if $H_{0}$ in (11) holds, then as $n \rightarrow \infty$

$$
H=\Upsilon^{-1} m \lambda_{m}^{4 \hat{d}_{A L P}}\left(\hat{\theta}_{2, A L P E}-\hat{\theta}_{2, A L P E}^{R}\right)^{2} \xrightarrow{d} \chi_{1}^{2}
$$

where

$$
\Upsilon=\frac{\pi^{2}}{6}\left[L_{1} \Omega\left(\hat{d}_{A L P}\right)^{-1} L_{1}^{\prime}-L_{2} \Omega_{2}\left(\hat{d}_{A L P}\right)^{-1} L_{2}^{\prime}\right]
$$

for $L_{1}=(0,0,1,0)$ and $L_{2}=(0,1)$.

The proof of this theorem relies on the consistency of $\hat{d}_{A L P}$ and the fact that $\sqrt{m} \lambda^{2 d_{0}}\left(\hat{\theta}_{2, A L P E}-\right.$ $\left.\hat{\theta}_{2, A L P E}^{R}\right) \xrightarrow{d} N(0, \Upsilon)$, which can be shown using the details in the proofs of Theorem 2 (for $\hat{\theta}_{2, A L P E}$ ) and of Theorem 3 in Arteche (2006) (for $\hat{\theta}_{2, A L P E}^{R}$ ) on the convergence of the respective Hessians and scores, extended to the nonstationary case (see the proof in Appendix A for details). This implies that the asymptotic covariance between $\hat{\theta}_{2, A L P E}^{R}$ and $\hat{\theta}_{2, A L P E}$ is equal to the variance of $\hat{\theta}_{2, A L P E}^{R}$, which is efficient in the class of augmented log periodogram regression estimators under uncorrelation of signal and noise. This structure is similar to other Hausman type tests.

The restricted $\hat{d}_{A L P}^{R}$ could have been used instead of $\hat{d}_{A L P}$ in the normalizing factor since it is also consistent under the null, but we use $\hat{d}_{A L P}$ due to its consistency under the null and the alternative. As before, $m^{-1} \lambda^{-4 \hat{d}_{A L P}} L_{1} \Omega\left(\hat{d}_{A L P}\right)^{-1} L_{1}^{\prime}$ and $m^{-1} \lambda^{-4 \hat{d}_{A L P}} L_{2} \Omega_{2}\left(\hat{d}_{A L P}\right)^{-1} L_{2}^{\prime}$ can be replaced by finite sample Hessian based approximations $L_{1} \Xi\left(\hat{d}_{A L P}, \hat{\theta}_{A L P}\right)^{-1} L_{1}^{\prime}$ and $L_{2} \Xi_{2}\left(\hat{d}_{A L P}, \hat{\theta}_{A L P}\right)^{-1} L_{2}^{\prime}$. In this case there is no need for the normalizing factor $m^{-1} \lambda^{-4 d_{0}}$ in the construction of the $H$ statistic.

## 6 FINITE SAMPLE PERFORMANCE

The finite sample performance of the proposed ALPE, which is denoted here by ALPEsin, is compared with the following estimators:

- The modified Gaussian semiparametric or local Whittle estimator of Hurvich et al. (2005) in their (P1) specification, i.e. ignoring the correlation between signal and noise. This estimator is denoted by MGSE.
- The modified Gaussian semiparametric or local Whittle estimator of Hurvich et al. (2005) in their (P2) specification, i.e. accounting only for the low frequency correlation and ignoring the high frequency correlation as if $\theta_{30}=0$. This estimator is denoted by MGSEcos.
- The modified Gaussian semiparametric estimator accounting for both the low and high frequency correlation as in formula (9). We denote this estimator by MGSEsin.
- The ALPE ignoring the correlation as suggested by Arteche (2006).
- The ALPE accounting only for the low frequency correlation as the MGSEcos (ALPEcos).

Under independence of signal and noise Arteche (2006) shows that the MGSE and the ALPE outperforms the estimators that do not account for the added noise (see also Hurvich and Ray, 2003). No results exist however on the correlated signal and noise case. Although this situation is partially covered in Hurvich et al. (2005) with the MGSEcos estimator, their Monte Carlo only considers MGSE in the independent case. The ALPEsin and MGSEsin are then new proposals and are expected to have a lower bias under general correlation structures. This assertion is confirmed by analyzing different situations. We first consider

$$
\begin{equation*}
z_{t}=\sigma_{y} y_{t}+u_{t} \tag{14}
\end{equation*}
$$

where $(1-L)^{d} y_{t}=w_{t}$ and $\sigma_{y}$ is chosen such that the long run noise to signal ratio is $n s r=2 \pi^{2}$, (similar results, available upon request, are obtained for $n s r=\pi^{2}$ ), which is close to the ratios considered by Arteche (2006), Deo and Hurvich (2001), Hurvich and Ray (2003) and Sun and Phillips (2003). We only show the results for $d=0.4$. Qualitatively similar conclusions derive for other values of the memory parameter. We have in particular analyzed also $d_{0}=0.2$ and $d_{0}=0.7$ (results available upon request) and the only difference is that, as expected, the larger $d_{0}$ the larger the bandwidth allowed in every estimator but the effects here discussed for $d=0.4$ remain unaltered. With this definition of $z_{t}$ two different scenarios are explored

Model 1: $u_{t}=\varepsilon_{t}$ and
Model 2: $u_{t}=\varepsilon_{t}+0.8 \varepsilon_{t-1}$
for

$$
\binom{\varepsilon_{t}}{w_{t}} \sim N I D\left(\binom{0}{0},\left(\begin{array}{ll}
1 & \rho \\
\rho & 1
\end{array}\right)\right)
$$

and $\rho=0,-0.8$.
Finally a non contemporaneous correlation example corresponding to an LMSV model is discussed. In this case the Gaussianity assumption of the added noise does not hold but we consider it relevant to analyze the applicability of the ALPE also in this context for its empirical interest. Its performance is compared with the local Whittle extensions whose asymptotic properties do not rely on the Gaussianity neither of the signal nor the noise.

Model 3: $z_{t}=\sigma_{y} y_{t-1}+u_{t}$ with $y_{t}$ defined as in Model 1 and $u_{t}=\log \varepsilon_{t}^{2}$ for $\varepsilon_{t}$ standard normal. For the sake of brevity we only show the results for the case of a correlation between
$w_{t}$ and $u_{t}$ equal to -0.8 . The results with null correlation (available upon request) are similar to those in Model 1. Note that correlation between $w_{t}$ and $\varepsilon_{t}$ is possible maintaining the null correlation between $w_{t}$ and $u_{t}$.

For Models 1, 2 and $3, \sigma_{y}^{2}=n s r^{-1},(1+0.8)^{2} n s r^{-1}$ and $n s r^{-1} \pi^{2} / 2$ respectively such that the signal to noise ratio is the same in all three models. In Model $1 \theta_{30}=0$ and the MGSEcos and ALPEcos are expected to perform better than the ALPEsin. However, in Models 2 and 3 with $\rho \neq 0$ there is a non null $\theta_{30}$ such that its omission increases the bias for large bandwidths and the ALPEsin and MGSEsin are expected to have a lower bias. The ALPE and MGSE are explicitly designed for $\rho=0$ and their behaviour when $\rho=-0.8$ is expected to be worse than the other estimators, at least in terms of bias.

A negative correlation raises a practical complication in the application of the different correlation corrected estimators. Whereas $1+\theta_{0}^{\prime} X_{j}\left(d_{0}\right)$ is always positive for a large enough sample size, in finite samples it can be negative, even when evaluated at the true set of parameter values, which prevents the logarithms from being taken in the objective functions. To circumvent this problem we truncate the argument in the logarithms by considering instead $\max \left(1+\theta^{\prime} X_{j}(d), 10^{-200}\right)$, which is asymptotically equivalent to $1+\theta^{\prime} X_{j}(d)$ for a large enough $n$. Note also that a negative $\rho$ implies that the real part of the cross spectral density diverges to $-\infty$ as $\lambda$ approaches the origin, greatly affecting the spectral behaviour of $z_{t}$ and the estimation of the parameters. We have also performed a similar analysis with a positive $\rho$ and the results (not reported but available upon request) show that the benefits of the correlation correction are not as evident as with a negative correlation; the bias decreases significantly but the variance inflation often gives rise to a higher mean square error.

The Monte Carlo consists of 1000 replications of series composed of 4096 observations. We choose such a large sample size to minimize the effect of the truncation of $1+\theta^{\prime} X_{j}(d)$ and because it is similar to the sample sizes of many of the financial time series which have formed the basis of several empirical applications on perturbed long memory as that analyzed in the next section. We analyze three different bandwidths, $m=n^{0.4}, n^{0.6}, n^{0.8}$. The latter increases at the same rate as the unfeasible ALPE optimal bandwidth but it can actually be far from this quantity due to the unknown multiplicative constant in $m_{A L P E}^{o p t}$. There does not exist however a feasible version of $m_{A L P E}^{o p t}$ and we do not pursue the issue here. Plug-in versions cannot be justified because they need an appropriate estimate of $G_{y}$, which is so
far not available. Criteria based on the minimization of an objective function may be a better choice but their performance is often unsatisfactory (Arteche, 2004). Other adaptive procedures such as that in Giraitis et al. (2000), which adapt to the spectral smoothness of $z_{t}$ could also be used, but any other bandwidth computed as the adaptive bandwidth times a constant would be equally (asymptotically) efficient. Recent results (Arteche and Orbe, 2009) using the bootstrap seem promising and can be easily extended to the ALPE based estimators but further research is required.

Table 2 shows the Monte Carlo biases and mean squared errors (MSE) in each of the different situations considered. For the minimization of the objective functions we have used the option nlminb in R , with the following restrictions: $0.01<d<0.9$, $-\exp (6)<\theta_{1}, \theta_{3}<$ $\exp (6)$ and $\exp (-20)<\theta_{2}<\exp (6)$.

When $\rho=0$ MGSE and ALPE tend to perform better and the inclusion of the terms accounting for correlation inevitably inflates the variance. The bias is however reduced if the noise is an $\mathrm{MA}(1)$ because the regressors for the correlation account indirectly for the weak dependence of the noise when a large bandwidth is used. When the signal and noise are correlated the ALPE and MGSE have lower MSE if a small bandwidth is used but the bias in both cases is quite large and tends to increase with $m$. The bias significantly decreases with the correlation correction such that for $m=n^{0.8}$ the ALPEsin tends to be the best option, not only in terms of bias but also in MSE.

The MGSE based estimators tend to perform poorly when correlation exists and a large bandwidth is used, even if the correlation is accounted for. This may be caused by the truncation of $1+\theta^{\prime} X_{j}(d)$ since it is at high frequencies that it is expected to have more impact and it plays its role in two terms of the contrast function in the MGSEcos and MGSEsin but only in one in the corresponding functions of the ALPEsin and ALPEcos. ${ }^{2}$

Finally, we analyze the performance of the different testing procedures for the hypothesis of no correlation between signal and noise. In the context considered in this paper the possible correlation is quite hard to detect because both series are unobservable. Moreover only the local spectral behaviour around the origin is restricted, thus permitting a great deal of flexibility. Table 3 shows the rejection frequencies of the Lagrange multiplier and

[^2]Table 2: Bias and MSE with $n s r=2 \pi^{2}$

|  | ALPE | ALPEcos | ALPEsin | MGSE | MGSEcos | MGSEsin |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | $u_{t}$ white noise$\rho=\mathbf{0}$ |  |  |  |  |  |
|  | $\mathbf{m}=\mathbf{n}^{0.4}$ |  |  |  |  |  |
| Bias | 0.050 | -0.127 | -0.152 | 0.000 | -0.130 | -0.121 |
| MSE | 0.042 | 0.093 | 0.094 | 0.036 | 0.095 | 0.091 |
|  | $\mathrm{m}=\mathrm{n}^{0.6}$ |  |  |  |  |  |
| Bias | 0.019 | -0.076 | -0.128 | -0.013 | -0.072 | -0.076 |
| MSE | 0.038 | 0.080 | 0.083 | 0.030 | 0.070 | 0.073 |
|  | $\mathbf{m}=\mathbf{n}^{0.8}$ |  |  |  |  |  |
| Bias | 0.021 | -0.026 | -0.062 | -0.001 | -0.023 | -0.037 |
| MSE | 0.024 | 0.057 | 0.074 | 0.016 | 0.033 | 0.042 |
|  | $\rho=-0.8$ |  |  |  |  |  |
|  | $\mathbf{m}=\mathbf{n}^{0.4}$ |  |  |  |  |  |
| Bias | 0.136 | -0.073 | -0.100 | 0.098 | -0.041 | -0.038 |
| MSE | 0.054 | 0.088 | 0.093 | 0.041 | 0.090 | 0.092 |
|  | $\mathrm{m}=\mathrm{n}^{0.6}$ |  |  |  |  |  |
| Bias | 0.165 | 0.008 | -0.070 | 0.166 | 0.174 | 0.162 |
| MSE | $\mathbf{m}=\mathbf{n}^{0.8}$ |  |  |  |  |  |
|  |  |  |  |  |  |  |
| Bias | 0.135 | -0.107 | 0.064 | 0.113 | 0.186 | 0.182 |
| MSE | 0.097 | 0.091 | 0.059 | 0.067 | 0.091 | 0.088 |
|  | $\begin{gathered} \mathbf{u}_{\mathbf{t}} \sim \mathbf{M A}(\mathbf{1}) \\ \rho=\mathbf{0} \end{gathered}$ |  |  |  |  |  |
|  | $\mathbf{m}=\mathbf{n}^{0.4}$ |  |  |  |  |  |
| Bias | 0.047 | -0.128 | -0.153 | 0.003 | -0.149 | -0.147 |
| MSE | 0.043 | 0.093 | 0.094 | 0.036 | 0.098 | 0.096 |
|  | $\mathrm{m}=\mathbf{n}^{0.6}$ |  |  |  |  |  |
| Bias | 0.017 | -0.076 | -0.135 | -0.015 | -0.092 | -0.112 |
| MSE | 0.036 | 0.080 | 0.084 | 0.031 | 0.076 | 0.082 |
|  | m $=\mathbf{n}^{\mathbf{0 . 8}}$ |  |  |  |  |  |
| Bias | -0.149 | -0.050 | -0.065 | -0.166 | -0.053 | -0.015 |
| MSE | 0.038 | 0.047 | 0.066 | 0.038 | 0.037 | 0.043 |
|  | $\rho=-0.8$ |  |  |  |  |  |
|  | $\mathbf{m}=\mathbf{n}^{0.4}$ |  |  |  |  |  |
| Bias | 0.129 | -0.069 | -0.103 | 0.096 | -0.051 | -0.051 |
| MSE | 0.053 | 0.088 | 0.091 | 0.042 | 0.092 | 0.093 |
|  | $\mathrm{m}=\mathrm{n}^{0.6}$ |  |  |  |  |  |
| Bias | 0.160 | 0.020 | -0.060 | 0.162 | 0.157 | 0.130 |
| MSE | 0.056 | 0.070 | 0.074 | 0.049 | 0.070 | 0.073 |
|  | m $\quad \mathbf{m} \mathbf{n}^{0.8}$ |  |  |  |  |  |
| Bias | 0.125 | 0.152 | -0.016 | 0.130 | 0.148 | 0.158 |
| MSE | 0.058 | 0.064 | 0.055 | 0.051 | 0.054 | 0.063 |
|  | $\begin{gathered} \mathbf{z}_{\mathbf{t}}=\sigma_{\mathbf{y}} \mathbf{y}_{\mathbf{t}-\mathbf{1}}+\log \chi_{\mathbf{1}}^{2} \\ \rho=-\mathbf{0 . 8} \end{gathered}$ |  |  |  |  |  |
|  | $\mathbf{m}=\mathbf{n}^{0.4}$ |  |  |  |  |  |
| Bias | 0.123 | -0.068 | -0.100 | 0.093 | -0.080 | -0.081 |
| MSE | 0.051 | 0.089 | 0.093 | 0.043 | 0.092 | 0.094 |
|  | $\mathrm{m}=\mathrm{n}^{\mathbf{0 . 6}}$ |  |  |  |  |  |
| Bias | 0.163 | 0.005 | -0.068 | 0.167 | 0.160 | 0.120 |
| MSE | 0.054 | 0.071 | 0.078 | 0.047 | 0.071 | 0.074 |
|  | $\mathrm{m}=\mathbf{n}^{0.8}$ |  |  |  |  |  |
| Bias | -0.358 | -0.369 | 0.116 | -0.164 | -0.124 | -0.098 |
| MSE | 0.148 | 0.146 | 0.046 | 0.112 | 0.131 | 0.114 |

Hausman test statistics for the null hypothesis of no correlation in (11). The nominal significance level is $5 \%$ (compared with the critical values of $\chi_{2}^{2}$ and $\chi_{1}^{2}$ distributions) and a bandwidth $m=n^{0.8}$ is used (we have found that this large bandwidth gives better results). Wald type tests are not considered because their performance is rather poor as mentioned in the previous section. The Lagrange multiplier has two important advantages over the Wald and Hausman type tests. First, no estimation of $\theta_{10}$ and $\theta_{30}$ is required. Second, since everything is calculated under the null of no correlation no truncation of the arguments of the logarithms is needed in the contrast functions.

To construct the test statistics a feasible approximation of $\Omega\left(d_{0}\right)$ and $\Omega_{2}\left(d_{0}\right)$ is needed. Two options can be considered: First a plug-in version replacing $d_{0}$ by the corresponding consistent estimate has the advantage of requiring only estimation of $d_{0}$, which gives rise to greater stability. Second we can use finite sample Hessian based approximations, as suggested by Sun and Phillips (2003) and Hurvich and Ray (2003), in two different forms: $J(d, \theta)+(H(d, \theta)-J(d, \theta)) I(H(d, \theta)>0)$ or just $J(d, \theta)$ as defined in the proof of Theorem 2 in Appendix A, with $d, \theta=\hat{d}_{A L P}^{R}, \hat{\theta}_{A L P}^{R}$ for the LM test and $d, \theta=\hat{d}_{A L P}, \hat{\theta}_{A L P}$ for the Hausman test statistic. In this case we need to estimate not only $d_{0}$ but also $\theta_{0}$ (only $\theta_{20}$ in the LM test), which makes the approximation less stable. However, when using either of them in the construction of the test statistics, the normalizing matrix $D_{n}^{*}$ in the LM and $m \lambda_{m}^{4 d_{0}}$ in the Hausman statistic are not needed due to the normalized convergence of both approximations to $\Omega\left(d_{0}\right)$ and $\Omega_{2}\left(d_{0}\right)$. Moreover we have found a better finite sample performance if either of these Hessian based approximation is used instead of the plug-in version and the rejection frequencies in Table 3 are obtained with both of them. The LM performs quite well if the added noise is white noise, particularly in the non-contemporaneous case. However, if the added noise shows some weak dependence, the LM testing procedure is not able to discriminate between the weak dependence of the noise and the correlation of signal and noise, since both arise as extra terms in the spectral density function. This situation could be corrected by approximating the spectral densities of the weak dependent innovations of signal and noise by local polynomials of finite orders, instead of constants. This would imply extending the nonlinear log periodogram regression with further elements to account for the weak dependence. By contrast, the Hausman test seems to be quite robust to weak dependence of the noise but is more conservative than the LM with low power in every situation.

Table 3: Rejection frequencies of $H_{0}: \theta_{10}=\theta_{30}=0(5 \%$ sig. level $)$ with $n s r=2 \pi^{2}$

| $\rho$ | $u_{t}$ white noise | $u_{t} \sim M A(1)$ | $z_{t}=\sigma_{y} y_{t-1}+\log \chi_{1}^{2}$ |
| :---: | :---: | :---: | :---: |
|  | LM |  |  |
| 0 | $0.095(0.096)$ | $0.340(0.341)$ | $0.076(0.077)$ |
| -0.8 | $0.889(0.889)$ | $0.554(0.554)$ | $0.999(0.999)$ |
|  | Hausman |  |  |
| 0 | $0.090(0.069)$ | $0.017(0.020)$ | $0.109(0.078)$ |
| -0.8 | $0.339(0.266)$ | $0.330(0.261)$ | $0.387(0.253)$ |

Rejection frequencies with $m=n^{0.8}$ using Hessian approximations (using $J$ between brackets).

## 7 EMPIRICAL EXAMPLE: S\&P500 REALIZED VOLATILITY

We analyze the persistence of the daily realized volatility (RV) for the S\&P500 future index. We construct the RV series using intraday transaction prices of futures contracts as traded on the Chicago Mercantile Exchange (CME) from 8:30AM to 3:15PM. As in Martens et al. (2009) the series is computed as the sum of squared intraday returns plus the squared overnight return between closing price and opening price next day, with a five minute sampling frequency. The sample period is from January 3, 1994 until May 29, 2009, extending the series analyzed by Martens et al. (2009), which runs until December 29, 2006. We also omit incomplete days and the large negative return for the period September 11-17, 2001. The series comprises a total of 3837 days. As shown by De Jong et al. (1998) transaction costs and lagged adjustment to information give rise to correlation between the underlying price and noise in stock price series that can transfer to the RV. We analyze here if such a correlation between the latent volatility and the noise persists in the S\&P500 RV and its effects on the estimation of the memory of the series.

Martens et al. (2009) use parametric techniques to get an estimate of the memory parameter around 0.5 . Table 4 shows the estimates for bandwidths $m=100,200$ and 300 using the six semiparametric estimation techniques considered in the Monte Carlo, together with the standard errors calculated by means of the matrix $J(d, \theta)$ as expalined in the previous section. ${ }^{3}$. There is a large discrepancy between the estimates accounting for the correlation and those ignoring it. Whereas the ALPE and the MGSE lies on the nonstationary region, the estimates accounting for the correlation shed some doubt even on

[^3]the strong persistence of the series. As shown in the Monte Carlo the correlation between signal and noise may induce a positive bias in the ALPE and MGSE and also in the ALPEcos and MGSEcos if a large bandwidth is used. We observe such a behaviour in Table 4 where estimates that account for correlation are similarly low with $m=100$ but the ALPEcos and MGSEcos increase significantly with the bandwidth whereas the ALPEsin and MGSEsin remain close to zero. The possible existence of correlation between signal and noise is corroborated by the LM and Hausman type tests statistics, which often give values larger than $10^{5}$, clearly supporting the existence of such correlation.

Table 4: Memory parameter estimates of S\&P500 RV and tests of correlation

| $m$ | ALPE | ALPEcos | ALPEsin | MGSE | MGSEcos | MGSEsin | $L M$ | $H$ |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 100 | 0.703 | 0.027 | 0.027 | 0.695 | 0.023 | 0.023 | $>10^{5}\left(>10^{5}\right)$ | $>10^{5}\left(>10^{5}\right)$ |
| s.e. | 0.129 | 0.295 | 1.237 | 0.101 | 0.272 | 1.116 |  |  |
| 200 | 0.626 | 0.289 | 0.107 | 0.609 | 0.281 | 0.082 | $3933.9(3933.9)$ | $>10^{5}\left(>10^{5}\right)$ |
| s.e. | 0.091 | 0.184 | 0.162 | 0.072 | 0.153 | 0.179 |  |  |
| 300 | 0.700 | 0.299 | 0.082 | 0.705 | 0.308 | 0.092 | $21.087(21.087)$ | $>10^{5}\left(>10^{5}\right)$ |
| s.e. | 0.085 | 0.136 | 0.218 | 0.064 | 0.116 | 0.211 |  |  |

$L M$ and $H$ are Lagrange Multiplier and Hausman test statistics for the null hypothesis of no correlation using Hessian approximations (using $J$ between brackets).

## APPENDIX A: PROOFS

Proof of Theorem 1: The method of proof in Sun and Phillips (2003) is used to avoid the flatness of $Q(d, \theta)$ as a function of $\theta$. Here we need the lemmas in Appendix B to correctly account for the correlation between signal and noise and also to avoid the linearization of the logarithm term in the contrast function, which, as explained in the text, would introduce a higher order bias. Write

$$
\begin{equation*}
Q(d, \theta)-Q\left(d_{0}, \theta_{0}\right)=\frac{1}{m} \sum_{j=1}^{m}\left(V_{j}^{\dagger}\right)^{2}+\frac{2}{m} \sum_{j=1}^{m} V_{j}^{\dagger}\left(U_{z j}+\varsigma_{j}\right) \tag{15}
\end{equation*}
$$

where

$$
\begin{aligned}
V_{j} & =V_{j}(d, \theta)=2\left(d-d_{0}\right) \log \lambda_{j}+\log \left(1+C_{j}^{0}\right)-\log \left(1+C_{j}\right) \\
C_{j} & =C_{j}(d, \theta)=\theta^{\prime} X_{j}(d) \\
C_{j}^{0} & =C_{j}\left(d_{0}, \theta_{0}\right)
\end{aligned}
$$

The consistency of $\hat{d}_{A L P}$ is established first. Since $\frac{1}{m} \sum_{j=1}^{m}\left(V_{j}^{\dagger}\right)^{2}=4\left(d-d_{0}\right)^{2}(1+o(1))$, where $o()$ holds uniformly in $\Delta \times \Theta$, we have to show that $\frac{2}{m} \sum_{j=1}^{m} V_{j}^{\dagger}\left(U_{z j}+\varsigma_{j}\right)=o_{p}(1)$
uniformly. First, by Lemmas 3 and 4 below

$$
\sup _{\Delta \times \Theta}\left|\sum_{j=1}^{m} V_{j}^{\dagger} U_{z j}\right|=O_{p}\left[\left(\frac{1}{\sqrt{m}}+\frac{\log ^{2} m}{m^{2\left(1-d_{0}\right)}}\right)\left[\sup _{\Delta}\left|d-d_{0}\right|+\sup _{\Delta} \lambda_{m}^{d}\right]\right]=o_{p}(1)
$$

Also by the definition of $\varsigma_{j}$ and by summation by parts

$$
\sup _{\Delta \times \Theta} \frac{1}{m}\left|\sum_{j=1}^{m} V_{j}^{\dagger} \varsigma_{j}\right|=O\left(\lambda_{m}^{\beta_{1}}\right)=o(1)
$$

which, together with the previous results, implies $\hat{d}_{A L P}-d_{0}=o_{p}(1)$.
Next, consider $d_{0}<3 / 4$. Since $Q\left(\hat{d}_{A L P}, \hat{\theta}_{A L P}\right)-Q\left(d_{0}, \theta_{0}\right) \leq 0$ then

$$
\begin{equation*}
\frac{1}{m} \sum{\widehat{V^{\dagger}}}_{j}^{2} \leq-\frac{2}{m} \sum{\widehat{V^{\dagger}}}_{j}\left(U_{z j}+\varsigma_{j}\right) \tag{16}
\end{equation*}
$$

where $\widehat{V}^{\dagger}{ }_{j}=V_{j}^{\dagger}\left(\hat{d}_{A L P}, \hat{\theta}_{A L P}\right)$. By Lemmas 3 and 4 and summation by parts the right hand side is bounded in probability by

$$
O_{p}\left(\frac{1}{\sqrt{m}}+\lambda_{m}^{\beta_{1}}\right)=o_{p}\left(\lambda_{m}^{2 \Delta_{1}}\right)
$$

where the second bound comes from the assumption $n^{2\left(d_{0}+1\right)(1+\delta)} m^{-2\left(d_{0}+1\right)(1+\delta)-1}=o(1)$. Since $m^{-1} \sum{\widehat{V^{\dagger}}}_{j}^{2}=4\left(\hat{d}_{A L P}-d_{0}\right)^{2}(1+o(1))+O\left(\lambda_{m}^{2 \Delta_{1}}\right)$ we have that $\hat{d}_{A L P}-d_{0}=O_{p}\left(\lambda_{m}^{\Delta_{1}}\right)=$ $o_{p}\left(\lambda_{m}^{\Delta_{1} / 2}\right)$. Consider now $d \in \Delta_{n}^{1}=\left\{d:\left|d-d_{0}\right|<\kappa \lambda_{m}^{\Delta_{1} / 2}\right\}$ for some generic constant $\kappa>0$ and $(d, \theta) \in \Delta_{n}^{1} \times \Theta$. Since by summation by parts $\frac{1}{m} \sum_{j=1}^{m} V_{j}^{\dagger} \varsigma_{j}$ equals

$$
\begin{align*}
& \frac{2\left(d-d_{0}\right)}{m} \sum\left(\log j-\frac{1}{m} \sum \log k\right) \varsigma_{j}+\frac{1}{m} \sum \log ^{\dagger}\left(\frac{1+C_{j}^{0}}{1+C_{j}}\right) \varsigma_{j} \\
= & O\left(\left|d-d_{0}\right| \lambda_{m}^{2}+\left|\theta_{10}-\theta_{1}\right| \lambda_{m}^{2+d_{0}}+\left|\theta_{20}-\theta_{2}\right| \lambda_{m}^{2+2 d_{0}}+\left|\theta_{30}-\theta_{3}\right| \lambda_{m}^{3+d_{0}}\right) \tag{17}
\end{align*}
$$

uniformly over $(d, \theta) \in \Delta_{n}^{1} \times \Theta$, we have that by (16), (17) and Corollary 2

$$
\frac{1}{m} \sum_{j=1}^{m}\left(\widehat{V}_{j}^{\dagger}\right)^{2} \leq O_{p}\left(\frac{1}{\sqrt{m}} \lambda_{m}^{\Delta_{1} / 2}\right)+O\left(\lambda_{m}^{\Delta_{1} / 2+\beta_{1}}\right)=o_{p}\left(\lambda_{m}^{a(1+\delta)}\right)
$$

for $a=1+d_{0}$, because $n^{2 a(1+\delta)} m^{-2 a(1+\delta)-1}=o(1)$. Then, using Lemma $5, \hat{d}_{A L P}-d_{0}=$ $o_{p}\left(\lambda_{m}^{a(1+\delta) / 2}\right), \hat{\theta}_{1, A L P}-\theta_{10}=o_{p}\left(\lambda_{m}^{a(1+\delta) / 2-d_{0}}\right)$ and $\hat{\theta}_{2, A L P}-\theta_{20}=O_{p}\left(\lambda_{m}^{b}\right)$ for $b=\max (0, a(1+$ $\left.\delta) / 2-2 d_{0}\right)$. The rest of the proof is made sequentially as in the proof of Theorem 2 in Sun and Phillips (2003) noting Corollary 2 and Lemma 5 below.

Proof of Theorem 2: The first order conditions are

$$
S\left(\hat{d}_{A L P}, \hat{\theta}_{A L P}\right)=0
$$

where, omitting the dependence on $(d, \theta)$ and $n$ for easy of notation, and denoting the vector $x_{j}=\left(x_{0 j}, x_{1 j}, x_{2 j}, x_{3 j}\right)^{\prime}$

$$
S(d, \theta)=\sum_{j=1}^{m} x_{j}^{\dagger} W_{j}
$$

with

$$
\begin{aligned}
x_{0 j}=x_{0 j}(d, \theta) & =\left(2 \log \lambda_{j}-\frac{C_{j}^{d}}{1+C_{j}}\right) \\
x_{i j}=x_{i j}(d, \theta) & =-\frac{C_{j}^{\theta_{i}}}{1+C_{j}}, \quad i=1,2,3 \\
W_{j}=W_{j}(d, \theta) & =\log I_{z j}+d\left(2 \log \lambda_{j}\right)-\log \left(1+C_{j}\right)
\end{aligned}
$$

for

$$
\begin{aligned}
& C_{j}^{d}=\frac{\partial C_{j}(d, \theta)}{\partial d}=\theta_{1} \lambda_{j}^{d}\left[A_{j}(d) \log \lambda_{j}+A_{j}^{d}(d)\right]+2 \theta_{2} \lambda_{j}^{2 d} \log \lambda_{j}+\theta_{3} \lambda_{j}^{1+d}\left[B_{j}(d) \log \lambda_{j}+B_{j}^{d}(d)\right] \\
& C_{j}^{\theta_{1}}=\frac{\partial C_{j}(d, \theta)}{\partial \theta_{1}}=A_{j}(d) \lambda_{j}^{d} \\
& C_{j}^{\theta_{2}}=\frac{\partial C_{j}(d, \theta)}{\partial \theta_{2}}=\lambda_{j}^{2 d} \\
& C_{j}^{\theta_{3}}=\frac{\partial C_{j}(d, \theta)}{\partial \theta_{3}}=B_{j}(d) \lambda_{j}^{1+d} \\
& \text { with } A_{j}^{d}(d)=-\left(\lambda_{j}-\pi\right) B_{j}(d) / 2 \text { and } B_{j}^{d}(d)=\left(\lambda_{j}-\pi\right) A_{j}(d) / 2 .
\end{aligned}
$$

The elements of the Hessian matrix, $H=H(d, \theta)$ are:

$$
\begin{aligned}
H_{1,1} & =\sum_{j=1}^{m}\left(x_{0 j}^{\dagger}\right)^{2}-\sum_{j=1}^{m}\left(\frac{C_{j}^{d d}}{1+C_{j}}-\frac{\left(C_{j}^{d}\right)^{2}}{\left(1+C_{j}\right)^{2}}\right)^{\dagger} W_{j} \\
H_{1,(i+1)} & =H_{(i+1), 1}=\sum_{j=1}^{m} x_{0 j}^{\dagger} x_{i j}^{\dagger}-\sum_{j=1}^{m}\left(\frac{C_{j}^{d \theta_{i}}}{1+C_{j}}-\frac{C_{j}^{d} C_{j}^{\theta_{i}}}{\left(1+C_{j}\right)^{2}}\right)^{\dagger} W_{j} \quad \text { for } i=1,2,3 \\
H_{(i+1),(k+1)} & =\sum_{j=1}^{m} x_{i j}^{\dagger} x_{k j}^{\dagger}+\sum_{j=1}^{m}\left(\frac{C_{j}^{\theta_{i}} C_{j}^{\theta_{k}}}{\left(1+C_{j}\right)^{2}}\right)^{\dagger} W_{j} \quad \text { for } i, k=1,2,3
\end{aligned}
$$

where

$$
\begin{aligned}
C_{j}^{d d} & =\theta_{1} A_{j}(d) \lambda_{j}^{d} \log ^{2} \lambda_{j}+2 \theta_{1} A_{j}^{d}(d) \lambda_{j}^{d} \log \lambda_{j}+\theta_{1} A_{j}^{d d}(d) \lambda_{j}^{d}+4 \theta_{2} \lambda_{j}^{2 d} \log ^{2} \lambda_{j} \\
& +\theta_{3} B_{j}(d) \lambda_{j}^{1+d} \log ^{2} \lambda_{j}+2 \theta_{3} B_{j}^{d}(d) \lambda_{j}^{1+d} \log \lambda_{j}+\theta_{3} B_{j}^{d d}(d) \lambda_{j}^{1+d} \\
C_{j}^{d \theta_{1}} & =A_{j}(d) \lambda_{j}^{d} \log \lambda_{j}+\lambda_{j}^{d} A_{j}^{d}(d) \\
C_{j}^{d \theta_{2}} & =2 \lambda_{j}^{2 d} \log \lambda_{j} \\
C_{j}^{d \theta_{3}} & =B_{j}(d) \lambda_{j}^{1+d} \log \lambda_{j}+\lambda_{j}^{1+d} B_{j}^{d}(d)
\end{aligned}
$$

for $A_{j}^{d d}(d)=-\left(\lambda_{j}-\pi\right)^{2} A_{j}(d) / 4$ and $B_{j}^{d d}(d)=-\left(\lambda_{j}-\pi\right)^{2} B_{j}(d) / 4$. Then

$$
\left(\hat{d}_{A L P E}, \hat{\theta}_{A L P}^{\prime}\right)^{\prime}-\left(d_{0}, \theta_{0}^{\prime}\right)^{\prime}=-H^{-1}(\bar{d}, \bar{\theta}) S\left(d_{0}, \theta_{0}\right)
$$

for $\left|\left(\bar{d}, \bar{\theta}^{\prime}\right)^{\prime}-\left(d_{0}, \theta_{0}^{\prime}\right)^{\prime}\right| \leq\left|\left(\hat{d}_{A L P E}, \hat{\theta}_{A L P}^{\prime}\right)^{\prime}-\left(d_{0}, \theta_{0}^{\prime}\right)^{\prime}\right|$. Considering the parameter set $\Delta_{n} \times$ $\Theta_{n}=\left\{(d, \theta):\left|\lambda_{m}^{-\left(1+d_{0}\right) / 2}\left(d-d_{0}\right)\right|<\epsilon,\left|\lambda_{m}^{-1 / 2}\left(\theta_{1}-\theta_{10}\right)\right|<\epsilon,\left|\lambda_{m}^{-\left(1-d_{0}\right) / 2}\left(\theta_{2}-\theta_{20}\right)\right|<\right.$ $\epsilon$ and $\left.\left|\theta_{3}-\theta_{30}\right|<\epsilon\right\}$ for an arbitrary small $\epsilon>0$, the asymptotic normality is proved by showing that:
a) $\sup _{(d, \theta) \in \Delta_{n} \times \Theta_{n}}\left\|D_{n}^{-1}[H(d, \theta)-J(d, \theta)] D_{n}^{-1}\right\|=o_{p}(1)$,
b) $\sup _{(d, \theta) \in \Delta_{n} \times \Theta_{n}}\left\|D_{n}^{-1}\left[J(d, \theta)-J\left(d_{0}, \theta_{0}\right)\right] D_{n}^{-1}\right\|=o_{p}(1)$,
c) $\quad D_{n}^{-1} J\left(d_{0}, \theta_{0}\right) D_{n}^{-1} \longrightarrow \Omega$
d) $\quad D_{n}^{-1} S\left(d_{0}, \theta_{0}\right) \xrightarrow{d} N\left(-b, \frac{\pi^{2}}{6} \Omega\right)$
where $J(d, \theta)$ is a $4 \times 4$ matrix with elements $[J(d, \theta)]_{(i+1),(k+1)}=\sum_{j=1}^{m} x_{i j}^{\dagger} x_{k j}^{\dagger}$, for $i, k=$ $0,1,2,3$. The ( 1,1 ) element of the left hand side of a) is

$$
-\sup _{\Delta_{n} \times \Theta_{n}} \sum_{j=1}^{m}\left(\frac{C_{j}^{d d}}{1+C_{j}}-\frac{\left(C_{j}^{d}\right)^{2}}{\left(1+C_{j}\right)^{2}}\right)^{\dagger}\left(V_{j}+U_{z j}+\varsigma_{j}\right)
$$

The term involving $V_{j}$ is bounded by

$$
O\left(\frac{1}{m} \sup _{\Delta_{n} \times \Theta_{n}} \sum_{j=1}^{m} \lambda_{j}^{d} \log ^{2} n\left[\left(d-d_{0}\right) \log n+\lambda_{j}^{d_{0}}\right]\right)=\left(\log ^{3} n \lambda_{m}^{2 d_{0}}\right)=o(1)
$$

under Assumption 4. Similarly the summand with $\varsigma_{j}$ is also $o(1)$ since $\varsigma_{j}=O\left(\lambda_{j}^{\beta_{1}}\right)$. Finally the term involving $U_{z j}$ is $o_{p}(1)$ using Lemma 3. The other elements are proved to be $o_{p}(1)$ similarly. b) and c) are easily proved noting the restrictions in the parameter space $\Delta_{n} \times \Theta_{n}$ and approximating sums by integrals.

Finally in order to prove d) write $D_{n}^{-1} S\left(d_{0}, \theta_{0}\right)=m^{-1 / 2} \sum N_{j}\left(U_{z j}+\varsigma_{j}\right)$ for

$$
N_{j}=\left(x_{0 j}^{\dagger}, \frac{\lambda_{m}^{-d_{0}}}{\cos \left(d_{0} \pi / 2\right)} x_{1 j}^{\dagger}, \lambda_{m}^{-2 d_{0}} x_{2 j}^{\dagger}, \frac{\lambda_{m}^{-\left(1+d_{0}\right)}}{\sin \left(d_{0} \pi / 2\right)} x_{3 j}^{\dagger}\right)^{\prime}
$$

It is easily shown that $m^{-1 / 2} \sum N_{j} \varsigma_{j} \rightarrow-b$ approximating sums by integrals. The proof is completed by showing that for any vector $v=\left(v_{1}, v_{2}, v_{3}, v_{4}\right)^{\prime}, m^{-1 / 2} \sum v^{\prime} N_{j} U_{z j} \xrightarrow{d}$ $N\left(0, v^{\prime} \Omega v \pi^{2} / 6\right)$. Divide now the sum in three parts

$$
\frac{1}{\sqrt{m}} \sum_{j=1}^{m} v^{\prime} N_{j} U_{z j}=T_{1}+T_{2}+T_{3}
$$

where, for $\beta=\max \left(2,\left(1-d_{0}\right)^{-1}\right)$ and $1>\alpha>\max \left\{1 / 2,1 / 4\left(1-d_{0}\right)\right\}$

$$
T_{1}=\frac{1}{\sqrt{m}} \sum_{j=1}^{\left[\log ^{\beta} m\right]} v^{\prime} N_{j} U_{z j}, \quad T_{2}=\frac{1}{\sqrt{m}} \sum_{j=1+\left[\log ^{\beta} m\right]}^{\left[m^{\alpha}\right]} v^{\prime} N_{j} U_{z j}, \quad T_{3}=\frac{1}{\sqrt{m}} \sum_{j=1+\left[m^{\alpha}\right]}^{m} v^{\prime} N_{j} U_{z j}
$$

Using Lemma 2 we can proceed as in Hurvich et al (1998) to show that $T_{1}$ and $T_{2}$ are both $o_{p}(1)$. Finally, since the elements in the vector $N_{j}$ satisfy the assumptions in Lemma 6 and $m^{-1} \sum_{j=1+\left[m^{\alpha}\right]}^{m}\left(v^{\prime} N_{j}\right)^{2}=v^{\prime} \Omega v(1+o(1))$ we get the desired result.

Proof of Theorem 4: Note that

$$
\begin{aligned}
& \hat{\theta}_{2, A L P}-\theta_{20}=-L_{1} H^{-1}(\bar{d}, \bar{\theta}) S\left(d_{0}, \theta_{0}\right) \\
& \hat{\theta}_{2, A L P}^{R}-\theta_{20}=-L_{2} G^{-1}(\bar{d}, \bar{\theta}) R\left(d_{0}, \theta_{20}\right)
\end{aligned}
$$

where $G()$ and $R()$ are the Hessian matrix and the score of the restricted ALPE in Arteche (2006, page 2124) respectively. Defining $D_{n 1}=D_{n 1}\left(d_{0}\right)=\sqrt{m} \operatorname{diag}\left(1, \lambda_{m}^{2 d_{0}}\right)$ we have that

$$
\sqrt{m} \lambda_{m}^{2 d_{0}}\left(\hat{\theta}_{2, A L P}-\hat{\theta}_{2, A L P}^{R}\right)=A_{n} B_{n}
$$

where $A_{n}=\left(-L_{1} D_{n} H^{-1}(\bar{d}, \bar{\theta}) D_{n}, L_{2} D_{n 1} G^{-1}(\bar{d}, \bar{\theta}) D_{n 1}\right)$ and $B_{n}=\left(D_{n}^{-1} S\left(d_{0}, \theta_{0}\right), D_{n 1}^{-1} R\left(d_{0}, \theta_{20}\right)\right)^{\prime}$. The desired result follows from the following convergences under the null hypothesis

$$
\begin{aligned}
& A_{n} \xrightarrow{p}\left(-L_{1} \Omega\left(d_{0}\right)^{-1}, L_{2} \Omega_{2}\left(d_{0}\right)^{-1}\right) \\
& B_{n} \xrightarrow{d} N\left(0, \frac{\pi^{2}}{6} \Omega_{4}\left(d_{0}\right)\right)
\end{aligned}
$$

for

$$
\Omega_{4}\left(d_{0}\right)=\left(\begin{array}{cc}
\Omega\left(d_{0}\right) & \Omega_{5}\left(d_{0}\right) \\
\Omega_{5}\left(d_{0}\right)^{\prime} & \Omega_{2}\left(d_{0}\right)
\end{array}\right)
$$

and $\Omega_{5}\left(d_{0}\right)$ containing the first and third columns of $\Omega\left(d_{0}\right)$. The convergence of $A_{n}$ is shown in a), b) and c) in the proof of Theorem 2 for the part related with $H$ and as in formula (A.5) in Arteche (2006) for the terms concerning $G$. The weak convergence of $B_{n}$ is shown in d) in the proof of Theorem 2 and as the last formula in the Appendix of Arteche (2006). Finally

$$
\left(-L_{1} \Omega\left(d_{0}\right)^{-1}, L_{2} \Omega_{2}\left(d_{0}\right)^{-1}\right) \Omega_{4}\left(d_{0}\right)\binom{-\Omega\left(d_{0}\right)^{-1} L_{1}^{\prime}}{\Omega_{2}\left(d_{0}\right)^{-1} L_{2}^{\prime}}=L_{1} \Omega\left(d_{0}\right)^{-1} L_{1}^{\prime}-L_{2} \Omega_{2}\left(d_{0}\right)^{-1} L_{2}^{\prime}
$$

for the form of the matrix $\Omega_{5}\left(d_{0}\right)$.

## APPENDIX B: TECHNICAL LEMMAS

Lemma 1 is a variant of Lemma 1 in Sun and Phillips (2003) for the stationary case and Theorem 1 in Velasco (1999) for $d_{0} \in[1 / 2,1)$ and the proof is thus omitted. Lemma

2 extends Lemma 2 in Sun and Phillips (2003) to the nonstationary case and the proof is similar noting the bounds in Lemma 1 to uniformly control for the errors in the approximation of the covariances between normalized discrete Fourier transforms. Lemma 3 is a more general version of Lemma 3 in Sun and Phillips (2003), which we found more convenient for the proofs of the theorems. Its proof is similar using Lemma 2 and it is thus omitted.

Lemma 1 Let $v_{j}=w_{z j}=w_{z j} / f_{z}^{1 / 2}\left(\lambda_{j}\right)$. Under assumptions 1-3, for any sequences of positive integers $j$ and $k$ such that $1 \leq k<j \leq m$ for $m / n \rightarrow 0$ as $n \rightarrow \infty$,
a) $E\left(v_{j} \bar{v}_{j}\right)=1+O\left(j^{-1} \log j+j^{2 d_{0}-2} \log j\right)$,
b) $E\left(v_{j} v_{j}\right)=O\left(j^{-1} \log j+j^{2 d_{0}-2} \log j\right)$,
c) $E\left(v_{j} \bar{v}_{k}\right)=O\left(k^{-1} \log j+(j k)^{d_{0}-1} \log k\right)$,
d) $E\left(v_{j} v_{k}\right)=O\left(k^{-1} \log j+(j k)^{d_{0}-1} \log k\right)$.

Lemma 2 Under assumptions 1-3 for $m / n \rightarrow 0$ and $\beta=\max \left(2,\left(1-d_{0}\right)^{-1}\right)$,
a) $\operatorname{Cov}\left(U_{z j}, U_{z k}\right)=O\left(k^{-2} \log ^{2} j+(j k)^{2 d_{0}-2} \log ^{2} k\right)$, uniformly for $\log ^{\beta} m \leq k<j \leq m$,
b) $\lim _{n} \sup _{1 \leq j \leq m} E U_{z j}^{2}<\infty$,
c) $E U_{z j}=O\left(j^{-1} \log j++j^{2 d_{0}-2} \log j\right)$, uniformly for $\log ^{\beta} m \leq k<j \leq m$,
d) $\operatorname{Var}\left(U_{z j}\right)=\pi^{2} / 6+O\left(j^{-1} \log j+j^{2 d_{0}-2} \log j\right)$, uniformly for $\log ^{\beta} m \leq k<j \leq m$,

Lemma 3 Let $\left\{c_{j}\left(d, \theta^{\prime}\right\}_{j=1}^{m}\right.$ be a sequence of functions such that for some finite $b>0$

$$
\begin{aligned}
\sup _{\Delta \times \Theta}\left|c_{j}-c_{j-1}\right| & =O\left(k_{1 m} j^{-1}\right) \quad \text { uniformly for } 2 \leq j \leq m \\
\sup _{\Delta \times \Theta}\left|c_{m}\right| & =O\left(k_{2 m}\right), \sup _{\Delta \times \Theta}\left|c_{j}\right|=O\left(\max \left(k_{1 m}, k_{2 m}\right) \frac{\sqrt{m}}{\log ^{b} m}\right)
\end{aligned}
$$

uniformly for $1 \leq j \leq m$. Then

$$
\sup _{\Delta \times \Theta}\left|\frac{1}{m} \sum_{j=1}^{m} c_{j} U_{z j}\right|=O_{p}\left[\max \left(k_{1 m}, k_{2 m}\right)\left(\frac{1}{\sqrt{m}}+\frac{\log ^{2} m}{m^{2\left(1-d_{0}\right)}}\right)\right]
$$

Lemma 4 For $(d, \theta) \in(\Delta \times \Theta)$ and $V_{j}$ defined in the proof of Theorem1,
a) $\left|V_{j}-V_{j-1}\right|=O\left(j^{-1}\left[\left|d-d_{0}\right|+\lambda_{j}^{d}+\lambda_{j}^{d_{0}}\right]\right)$, uniformly for $2 \leq j \leq m$,
b) $\left|V_{j}^{\dagger}\right|=O\left(\left|d-d_{0}\right| \log m+\lambda_{m}^{d}+\lambda_{m}^{d_{0}}\right)$, uniformly for $1 \leq j<m$,
c) $\left.\left|V_{m}^{\dagger}\right|=O\left(\left|d-d_{0}\right|+\lambda_{m}^{d}+\lambda_{m}^{d_{0}}\right]\right)$,
uniformly in $\Delta \times \Theta$ and if $d \in \Delta_{n}^{\prime}=\left\{d:\left|d-d_{0}\right|<\kappa \lambda_{m}^{v}\right\}$ for some finite constant $\kappa$ and $v>0$ arbitrary small,
d) $\left|V_{j}-V_{j-1}\right|=O\left(j^{-1}\left[\left|d-d_{0}\right|+\left|\theta_{1}-\theta_{10}\right| \lambda_{m}^{d_{0}}+\left|\theta_{2}-\theta_{20}\right| \lambda_{m}^{2 d_{0}}+\left|\theta_{3}-\theta_{30}\right| \lambda_{m}^{1+d_{0}}\right)\right.$, uniformly for $2 \leq j \leq m$,
e) $\left|V_{j}^{\dagger}\right|=O\left(\left|d-d_{0}\right| \log m+\left|\theta_{1}-\theta_{10}\right| \lambda_{m}^{d_{0}}+\left|\theta_{2}-\theta_{20}\right| \lambda_{m}^{2 d_{0}}+\left|\theta_{3}-\theta_{30}\right| \lambda_{m}^{1+d_{0}}\right)$, uniformly for $1 \leq j<m$,
f) $\left|V_{m}^{\dagger}\right|=O\left(\left|d-d_{0}\right|+\left|\theta_{1}-\theta_{10}\right| \lambda_{m}^{d_{0}}+\left|\theta_{2}-\theta_{20}\right| \lambda_{m}^{2 d_{0}}+\left|\theta_{3}-\theta_{30}\right| \lambda_{m}^{1+d_{0}}\right)$,
uniformly in $\Delta_{n}^{\prime} \times \Theta$.

Proof: a) By the definition of $V_{j}$

$$
\begin{equation*}
\left|V_{j}-V_{j-1}\right|=\left|2\left(d-d_{0}\right) \log \left(\frac{j}{j-1}\right)+\log \left(\frac{1+C_{j}^{0}}{1+C_{j-1}^{0}}\right)-\log \left(\frac{1+C_{j}}{1+C_{j-1}}\right)\right| \tag{18}
\end{equation*}
$$

The first term of the right hand side of (18) is $O\left(\left|d-d_{0}\right| j^{-1}\right)$ where hereafter the $O()$ are uniform over $2 \leq j<m$ and $\Delta \times \Theta$ or $\Delta_{n}^{\prime} \times \Theta$ where appropriate. Now

$$
\log \left(\frac{1+C_{j}}{1+C_{j-1}}\right)=\log \left(1+\frac{C_{j}-C_{j-1}}{1+C_{j-1}}\right)=O\left(\left|C_{j}-C_{j-1}\right|\right)=O\left(\frac{1}{j} \lambda_{j}^{d}\right)
$$

using the following bounds

$$
\begin{aligned}
\left|\lambda_{j}^{\alpha}-\lambda_{j-1}^{\alpha}\right| & =\lambda_{j}^{\alpha}\left|1-\left(\frac{j-1}{j}\right)^{\alpha}\right|=O\left(\frac{1}{j} \lambda_{j}^{\alpha}\right) \quad \text { for any } \alpha \neq 0 \\
\cos \left[d\left(\frac{\lambda_{j}}{2}+\frac{\pi}{2}\right)\right] & =\cos \left[d\left(\frac{\lambda_{j-1}}{2}+\frac{\pi}{2}\right)\right]+O\left(\frac{1}{n}\right)
\end{aligned}
$$

and similarly for the sine term.
d) For the bound in d) note that

$$
\begin{equation*}
\left|V_{j}-V_{j-1}\right| \leq\left|2\left(d-d_{0}\right) \log \left(\frac{j}{j-1}\right)\right|+\left|\log \left(\frac{\left(1+C_{j}^{0}\right)\left(1+C_{j-1}\right)}{\left(1+C_{j-1}^{0}\right)\left(1+C_{j}\right)}\right)\right| \tag{19}
\end{equation*}
$$

where the first term on the right hand side is $O\left(j^{-1}\left|d-d_{0}\right|\right)$. The second term is the absolute value of

$$
\log \left(1+\frac{C_{j}^{0}-C_{j}-\left(C_{j-1}^{0}-C_{j-1}\right)+C_{j}^{0} C_{j-1}-C_{j-1}^{0} C_{j}}{\left(1+C_{j-1}^{0}\right)\left(1+C_{j}\right)}\right)
$$

Now $C_{j}^{0}-C_{j}-\left(C_{j-1}^{0}-C_{j-1}\right)$ is equal to

$$
\begin{align*}
& \left(\theta_{10}-\theta_{1}\right)\left[\lambda_{j}^{d_{0}} A_{j}\left(d_{0}\right)-\lambda_{j-1}^{d_{0}} A_{j-1}\left(d_{0}\right)\right]  \tag{20}\\
+ & \theta_{1}\left[A_{j}\left(d_{0}\right)\left(\lambda_{j}^{d_{0}}-\lambda_{j}^{d}\right)-A_{j-1}\left(d_{0}\right)\left(\lambda_{j-1}^{d_{0}}-\lambda_{j-1}^{d}\right)\right]  \tag{21}\\
+ & \theta_{1}\left[\lambda_{j}^{d}\left(A_{j}\left(d_{0}\right)-A_{j}(d)\right)-\lambda_{j-1}^{d}\left(A_{j-1}\left(d_{0}\right)-A_{j-1}(d)\right)\right]  \tag{22}\\
+ & \left(\theta_{20}-\theta_{2}\right)\left(\lambda_{j}^{2 d_{0}}-\lambda_{j-1}^{2 d_{0}}\right)  \tag{23}\\
+ & \theta_{2}\left[\lambda_{j}^{2 d_{0}}-\lambda_{j-1}^{2 d_{0}}-\left(\lambda_{j}^{2 d}-\lambda_{j-1}^{2 d}\right)\right]  \tag{24}\\
+ & \left(\theta_{30}-\theta_{3}\right)\left[\lambda_{j}^{1+d_{0}} B_{j}\left(d_{0}\right)-\lambda_{j-1}^{1+d_{0}} B_{j-1}\left(d_{0}\right)\right]  \tag{25}\\
+ & \theta_{3}\left[B_{j}\left(d_{0}\right)\left(\lambda_{j}^{1+d_{0}}-\lambda_{j}^{1+d}\right)-B_{j-1}\left(d_{0}\right)\left(\lambda_{j-1}^{1+d_{0}}-\lambda_{j-1}^{1+d}\right)\right]  \tag{26}\\
+ & \theta_{3}\left[\lambda_{j}^{1+d}\left(B_{j}\left(d_{0}\right)-B_{j}(d)\right)-\lambda_{j-1}^{1+d}\left(B_{j-1}\left(d_{0}\right)-B_{j-1}(d)\right)\right] \tag{27}
\end{align*}
$$

The expression in (20) is $\theta_{10}-\theta_{1}$ times

$$
\left(\lambda_{j}^{d_{0}}-\lambda_{j-1}^{d_{0}}\right) A_{j}\left(d_{0}\right)+\lambda_{j-1}^{d_{0}}\left(A_{j}\left(d_{0}\right)-A_{j-1}\left(d_{0}\right)\right)=O\left(\frac{1}{j} \lambda_{j}^{d_{0}}\right)+O\left(\frac{1}{n} \lambda_{j-1}^{d_{0}}\right)=O\left(\frac{1}{j} \lambda_{j}^{d_{0}}\right)
$$

Now, apart from the constant $\theta_{1},(21)$ is

$$
\begin{aligned}
& \left(\lambda_{j}^{d_{0}}-\lambda_{j}^{d}\right)\left[A_{j}\left(d_{0}\right)-A_{j-1}\left(d_{0}\right)\right]+A_{j-1}\left(d_{0}\right)\left[\lambda_{j}^{d_{0}}-\lambda_{j}^{d}-\left(\lambda_{j-1}^{d_{0}}-\lambda_{j-1}^{d}\right)\right] \\
= & O\left(\left|d-d_{0}\right| \lambda_{j}^{d_{0}} \log \lambda_{j} \frac{1}{n}\right)+O\left(\left|d-d_{0}\right| \lambda_{j}^{d_{0}} \log \lambda_{j} \frac{1}{j}\right) \\
= & O\left(\left|d-d_{0}\right| \lambda_{j}^{d_{0}} \log \lambda_{j} \frac{1}{j}\right)
\end{aligned}
$$

Similarly, (22) is, apart from $\theta_{1}$

$$
\begin{aligned}
& \left(\lambda_{j}^{d}-\lambda_{j-1}^{d}\right)\left[A_{j}\left(d_{0}\right)-A_{j}(d)\right]+\lambda_{j-1}^{d}\left[A_{j}\left(d_{0}\right)-A_{j-1}\left(d_{0}\right)-\left(A_{j}(d)-A_{j-1}(d)\right)\right] \\
= & O\left(\frac{1}{j} \lambda_{j}^{d}\left|d-d_{0}\right|\right)+\lambda_{j-1}^{d}\left[d_{0} A_{j-1}^{\prime}\left(d_{0}\right)-d A_{j-1}^{\prime}(d)+O\left(\frac{1}{n^{2}}\right)\right] \\
= & O\left(\frac{1}{j} \lambda_{j}^{d}\left|d-d_{0}\right|\right)+\lambda_{j-1}^{d}\left[O\left(\frac{1}{n^{2}}\right)+d_{0}\left[A_{j-1}^{\prime}\left(d_{0}\right)-A_{j-1}^{\prime}(d)\right]+\left(d_{0}-d\right) A_{j-1}^{\prime}(d)\right] \\
= & O\left(\frac{1}{j} \lambda_{j}^{d}\left|d-d_{0}\right|\right)
\end{aligned}
$$

for $A_{j-1}^{\prime}(d)=-\sin \left[d\left(\lambda_{j-1} / 2-\pi / 2\right)\right] \pi / n$. We get similarly the bounds for (25), (26) and (27) with $\lambda_{j}^{1+d}$ instead of $\lambda_{j}^{d}$.

Also (23) is $O\left(\left|\theta_{20}-\theta_{2}\right| j^{-1} \lambda_{j}^{2 d_{0}}\right)$ and, apart from the constant $\theta_{2},(24)$ is equal to

$$
\begin{aligned}
& \lambda_{j}^{2 d_{0}}\left[\left(1-\frac{1}{j}\right)^{2 d}-\left(1-\frac{1}{j}\right)^{2 d_{0}}\right]+\left[1-\left(1-\frac{1}{j}\right)^{2 d}\right]\left(\lambda_{j}^{2 d_{0}}-\lambda_{j}^{2 d}\right) \\
= & O\left(\lambda_{j}^{2 d_{0}}\left|d-d_{0}\right| \frac{1}{j}\right)+\left(\lambda_{j}^{2 d_{0}}\left|d-d_{0}\right| \frac{1}{j} \log \lambda_{j}\right) \\
= & O\left(\left|d-d_{0}\right| \frac{1}{j}\right)
\end{aligned}
$$

Finally $C_{j}^{0} C_{j-1}-C_{j-1}^{0} C_{j}$ is equal to

$$
\begin{align*}
& {\left[C_{j}^{0}-C_{j}-\left(C_{j-1}^{0}-C_{j-1}\right)\right] C_{j-1}^{0}+\left(C_{j}^{0}-C_{j-1}^{0}\right)\left(C_{j-1}-C_{j-1}^{0}\right) }  \tag{28}\\
= & O\left(C_{j-1}^{0} \frac{1}{j}\left[\left|d-d_{0}\right|+\left|\theta_{1}-\theta_{10}\right| \lambda_{m}^{d_{0}}+\left|\theta_{2}-\theta_{20}\right| \lambda_{m}^{2 d_{0}}+\left|\theta_{3}-\theta_{30}\right| \lambda_{m}^{1+d_{0}}\right)\right.
\end{align*}
$$

from the proof of the bounds in (20)-(27) and because

$$
\begin{aligned}
C_{j}^{0}-C_{j-1}^{0} & =\theta_{10} \lambda_{j}^{d_{0}}\left[A_{j}\left(d_{0}\right)-A_{j-1}\left(d_{0}\right)\right] \\
& +\theta_{10} A_{j-1}\left(d_{0}\right)\left(\lambda_{j}^{d_{0}}-\lambda_{j-1}^{d_{0}}\right)+\theta_{20}\left(\lambda_{j}^{2 d_{0}}-\lambda_{j-1}^{2 d_{0}}\right) \\
& +\theta_{30} \lambda_{j}^{1+d_{0}}\left[B_{j}\left(d_{0}\right)-B_{j-1}\left(d_{0}\right)\right] \\
& +\theta_{30} B_{j-1}\left(d_{0}\right)\left(\lambda_{j}^{1+d_{0}}-\lambda_{j-1}^{1+d_{0}}\right) \\
& =O\left(\frac{1}{n} \lambda_{j}^{d_{0}}\right)+O\left(\frac{1}{j} \lambda_{j}^{d_{0}}\right)+O\left(\frac{1}{j} \lambda_{j}^{2 d_{0}}\right)=O\left(\frac{1}{j} \lambda_{j}^{d_{0}}\right) \\
C_{j-1}^{0}-C_{j-1}= & \left(\theta_{10}-\theta_{1}\right) \lambda_{j-1}^{d_{0}} A_{j-1}\left(d_{0}\right)+\theta_{1} A_{j-1}\left(d_{0}\right)\left(\lambda_{j-1}^{d_{0}}-\lambda_{j-1}^{d}\right) \\
& +\theta_{1} \lambda_{j-1}^{d}\left[A_{j-1}\left(d_{0}\right)-A_{j-1}(d)\right]+\left(\theta_{20}-\theta_{2}\right) \lambda_{j-1}^{2 d_{0}}+\theta_{2}\left(\lambda_{j-1}^{2 d_{0}}-\lambda_{j-1}^{d}\right) \\
+ & \left(\theta_{30}-\theta_{3}\right) \lambda_{j-1}^{1+d_{0}} B_{j-1}\left(d_{0}\right)+\theta_{3} B_{j-1}\left(d_{0}\right)\left(\lambda_{j-1}^{1+d_{0}}-\lambda_{j-1}^{1+d}\right) \\
+ & \theta_{3} \lambda_{j-1}^{1+d}\left[B_{j-1}\left(d_{0}\right)-B_{j-1}(d)\right] \\
= & O\left(\left|\theta_{10}-\theta_{1}\right| \lambda_{m}^{d_{0}}+\lambda_{m}^{d_{0}}\left|d-d_{0}\right| \log \lambda_{m}+\left|\theta_{20}-\theta_{2}\right| \lambda_{m}^{2 d_{0}}+\left|\theta_{30}-\theta_{3}\right| \lambda_{m}^{1+d_{0}}\right)
\end{aligned}
$$

b) The bound for $V_{j}^{\dagger}$ in b) uses the following relations

$$
\begin{aligned}
\log \lambda_{j}-\frac{1}{m} \sum \log \lambda_{k} & =\log j-\frac{1}{m} \sum \log k=O(\log m) \\
\log \left(\frac{1+C_{j}^{0}}{1+C_{j}}\right) & =\log \left(1+\frac{C_{j}^{0}-C_{j}}{1+C_{j}}\right)=O\left(\left|C_{j}^{0}-C_{j}\right|\right)=O\left(\lambda_{m}^{d}+\lambda_{m}^{d_{0}}\right)
\end{aligned}
$$

e) For $d \in \Delta_{n}^{\prime}$ we get e) using also the following bounds

$$
\begin{aligned}
\lambda_{j}^{d} & =\lambda_{j}^{d_{0}}(1+o(1)) \quad, \quad \lambda_{j}^{2 d}=\lambda_{j}^{2 d_{0}}(1+o(1)) \\
A_{j}(d) & =A_{j}\left(d_{0}\right)+O\left(\left|d-d_{0}\right|\right) \quad, \quad B_{j}(d)=B_{j}\left(d_{0}\right)+O\left(\left|d-d_{0}\right|\right)
\end{aligned}
$$

The bounds of c) and f) are equally obtained noting that $\log \lambda_{m}-m^{-1} \sum \log \lambda_{k}=$ $\log m-m^{-1} \sum \log k=1+o(1)$.

Corollary 2 By Lemmas 3 and $4 \sup _{\Delta_{n}^{\prime} \times \Theta}\left|m^{-1} \sum_{j=1}^{m} V_{j}^{\dagger} U_{z j}\right|$ is bounded in probability by

$$
O_{p}\left[\left(\frac{1}{\sqrt{m}}+\frac{\log ^{2} m}{m^{2\left(1-d_{0}\right)}}\right) \sup _{\Delta_{n}^{\prime} \times \Theta}\left[\left|d-d_{0}\right|+\left|\theta_{10}-\theta_{1}\right| \lambda_{m}^{d_{0}}+\left|\theta_{20}-\theta_{2}\right| \lambda_{m}^{2 d_{0}}+\left|\theta_{30}-\theta_{3}\right| \lambda_{m}^{1+d_{0}}\right]\right]
$$

Lemma 5 For $(d, \theta) \in \Delta_{n}^{\prime} \times \Theta$

$$
\begin{aligned}
\frac{1}{m} \sum_{j=1}^{m}\left(V_{j}^{\dagger}\right)^{2} & =4\left(d-d_{0}\right)^{2}(1+o(1)) \\
& +\left(\theta_{10}-\theta_{1}\right)^{2} \cos ^{2}\left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{2 d_{0}} \frac{d_{0}^{2}}{\left(2 d_{0}+1\right)\left(1+d_{0}\right)^{2}}(1+o(1)) \\
& +\left(\theta_{20}-\theta_{2}\right)^{2} \lambda_{m}^{4 d_{0}} \frac{4 d_{0}^{2}}{\left(4 d_{0}+1\right)\left(1+2 d_{0}\right)^{2}}(1+o(1)) \\
& +\left(\theta_{30}-\theta_{3}\right)^{2} \sin ^{2}\left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{2\left(d_{0}+1\right)} \frac{\left(1+d_{0}\right)^{2}}{\left(2 d_{0}+3\right)\left(2+d_{0}\right)^{2}}(1+o(1)) \\
& +2\left(\theta_{10}-\theta_{1}\right)\left(\theta_{20}-\theta_{2}\right) \cos \left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{3 d_{0}} \frac{2 d_{0}^{2}}{\left(3 d_{0}+1\right)\left(d_{0}+1\right)\left(1+2 d_{0}\right)}(1+o(1)) \\
& -2\left(\theta_{10}-\theta_{1}\right)\left(\theta_{30}-\theta_{3}\right) \cos \left(\frac{d_{0} \pi}{2}\right) \sin \left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{1+2 d_{0}} \frac{d_{0}}{2\left(d_{0}+1\right)\left(d_{0}+2\right)}(1+o(1)) \\
& -2\left(\theta_{20}-\theta_{2}\right)\left(\theta_{30}-\theta_{3}\right) \sin \left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{1+3 d_{0}} \frac{2 d_{0}\left(1+d_{0}\right)}{\left(3 d_{0}+2\right)\left(2 d_{0}+1\right)\left(d_{0}+2\right)}(1+o(1)) \\
& +4\left(d-d_{0}\right)\left(\theta_{10}-\theta_{1}\right) \cos \left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{d_{0}} \frac{d_{0}}{\left(d_{0}+1\right)^{2}}(1+o(1)) \\
& +4\left(d-d_{0}\right)\left(\theta_{20}-\theta_{2}\right) \lambda_{m}^{2 d_{0}} \frac{2 d_{0}}{\left(2 d_{0}+1\right)^{2}}(1+o(1)) \\
& +4\left(d-d_{0}\right)\left(\theta_{30}-\theta_{3}\right) \lambda_{m}^{1+d_{0}} \frac{1+d_{0}}{\left(d_{0}+2\right)^{2}}(1+o(1))
\end{aligned}
$$

where the $o(1)$ terms are uniform over $(d, \theta) \in \Delta_{n}^{\prime} \times \Theta$.

## Proof:

$$
\begin{align*}
\frac{1}{m} \sum_{j=1}^{m}\left(V_{j}^{\dagger}\right)^{2} & =4\left(d-d_{0}\right)^{2} \frac{1}{m} \sum\left(\log ^{\dagger} \lambda_{j}\right)^{2}  \tag{29}\\
& +\frac{1}{m} \sum\left[\log ^{\dagger}\left(\frac{1+C_{j}^{0}}{1+C_{j}}\right)\right]^{2}  \tag{30}\\
& +\frac{4\left(d-d_{0}\right)}{m} \sum \log ^{\dagger} \lambda_{j} \log \left(\frac{1+C_{j}^{0}}{1+C_{j}}\right) \tag{31}
\end{align*}
$$

The right hand side of $(29)$ is $4\left(d-d_{0}\right)^{2}\left(1+O\left(m^{-1} \log ^{2} m\right)\right)$ where hereafter the $o()$ and
$O()$ terms hold uniformly over $(d, \theta) \in \Delta_{n}^{\prime} \times \Theta$. Now

$$
\log \left(\frac{1+C_{j}^{0}}{1+C_{j}}\right)=\log \left(1+\frac{C_{j}^{0}-C_{j}}{1+C_{j}}\right)=C_{j}^{0}-C_{j}-\frac{C_{j}\left(C_{j}^{0}-C_{j}\right)}{1+C_{j}}+O\left[\left(C_{j}^{0}-C_{j}\right)^{2}\right]
$$

and $C_{j}^{0}-C_{j}$ is equal to

$$
\left(\theta_{10}-\theta_{1}\right) A_{j}\left(d_{0}\right) \lambda_{j}^{d_{0}}+\left(\theta_{20}-\theta_{2}\right) \lambda_{j}^{2 d_{0}}+\left(\theta_{30}-\theta_{3}\right) B_{j}\left(d_{0}\right) \lambda_{j}^{1+d_{0}}+O\left(\left|d-d_{0}\right| \lambda_{j}^{d_{0}} \log \lambda_{j}\right)
$$

since $\lambda_{j}^{d}=\lambda_{j}^{d_{0}}+O\left(\left|d-d_{0}\right| \lambda_{j}^{d_{0}} \log \lambda_{j}\right)$ such that

$$
\begin{aligned}
\log \left(\frac{1+C_{j}^{0}}{1+C_{j}}\right) & =\left(\theta_{10}-\theta_{1}\right) A_{j}\left(d_{0}\right) \lambda_{j}^{d_{0}}+\left(\theta_{20}-\theta_{2}\right) \lambda_{j}^{2 d_{0}}+\left(\theta_{30}-\theta_{3}\right) B_{j}\left(d_{0}\right) \lambda_{j}^{1+d_{0}} \\
& +O\left(\left|d-d_{0}\right| \lambda_{j}^{d_{0}} \log \lambda_{j}+\left|\theta_{10}-\theta_{1}\right| \lambda_{j}^{2 d_{0}}+\left|\theta_{20}-\theta_{2}\right| \lambda_{j}^{3 d_{0}}+\left|\theta_{30}-\theta_{3}\right| \lambda_{j}^{1+2 d_{0}}\right)
\end{aligned}
$$

Then

$$
\begin{aligned}
\frac{1}{m} \sum \log ^{2}\left(\frac{1+C_{j}^{0}}{1+C_{j}}\right) & =\left(\theta_{10}-\theta_{1}\right)^{2} \cos ^{2}\left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{2 d_{0}} \frac{1}{2 d_{0}+1}(1+o(1)) \\
& +\left(\theta_{20}-\theta_{2}\right)^{2} \lambda_{m}^{4 d_{0}} \frac{1}{4 d_{0}+1}(1+o(1)) \\
& +\left(\theta_{30}-\theta_{3}\right)^{2} \sin ^{2}\left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{2\left(1+d_{0}\right)} \frac{1}{2 d_{0}+3}(1+o(1)) \\
& +2\left(\theta_{10}-\theta_{1}\right)\left(\theta_{20}-\theta_{2}\right) \cos \left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{3 d_{0}} \frac{1}{3 d_{0}+1}(1+o(1)) \\
& -2\left(\theta_{10}-\theta_{1}\right)\left(\theta_{30}-\theta_{3}\right) \cos \left(\frac{d_{0} \pi}{2}\right) \sin \left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{1+2 d_{0}} \frac{1}{2 d_{0}+2}(1+o(1)) \\
& -2\left(\theta_{20}-\theta_{2}\right)\left(\theta_{30}-\theta_{3}\right) \sin \left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{1+3 d_{0}} \frac{1}{3 d_{0}+2}(1+o(1)) \\
& +O\left(\lambda_{m}^{2 d_{0}}\left(d-d_{0}\right)^{2} \log ^{2} \lambda_{m}+\left|d-d_{0}\right|\left|\theta_{10}-\theta_{1}\right| \lambda_{m}^{2 d_{0}} \log \lambda_{m}\right. \\
& \left.+\left|d-d_{0}\right|\left|\theta_{20}-\theta_{2}\right| \lambda_{m}^{3 d_{0}} \log \lambda_{m}+\left|d-d_{0}\right|\left|\theta_{30}-\theta_{3}\right| \lambda_{m}^{1+2 d_{0}} \log \lambda_{m}\right)
\end{aligned}
$$

$$
\begin{aligned}
{\left[\frac{1}{m} \sum \log \left(\frac{1+C_{j}^{0}}{1+C_{j}}\right)\right]^{2} } & =\left(\theta_{10}-\theta_{1}\right)^{2} \cos ^{2}\left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{2 d_{0}} \frac{1}{\left(d_{0}+1\right)^{2}}(1+o(1)) \\
& +\left(\theta_{20}-\theta_{2}\right)^{2} \lambda_{m}^{4 d_{0}} \frac{1}{\left(2 d_{0}+1\right)^{2}}(1+o(1)) \\
& +\left(\theta_{30}-\theta_{3}\right)^{2} \sin ^{2}\left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{2+2 d_{0}} \frac{1}{\left(d_{0}+2\right)^{2}}(1+o(1)) \\
& +2\left(\theta_{10}-\theta_{1}\right)\left(\theta_{20}-\theta_{2}\right) \cos \left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{3 d_{0}} \frac{1}{\left(1+d_{0}\right)\left(2 d_{0}+1\right)}(1+o(1)) \\
& -2\left(\theta_{10}-\theta_{1}\right)\left(\theta_{30}-\theta_{3}\right) \cos \left(\frac{d_{0} \pi}{2}\right) \sin \left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{1+2 d_{0}} \frac{(1+o(1))}{\left(1+d_{0}\right)\left(d_{0}+2\right)} \\
& -2\left(\theta_{20}-\theta_{2}\right)\left(\theta_{30}-\theta_{3}\right) \sin \left(\frac{d_{0} \pi}{2}\right) \lambda_{m}^{1+3 d_{0}} \frac{1}{\left(1+2 d_{0}\right)\left(d_{0}+2\right)}(1+o(1)) \\
& +O\left(\lambda_{m}^{2 d_{0}}\left(d-d_{0}\right)^{2} \log ^{2} \lambda_{m}+\left|d-d_{0}\right|\left|\theta_{10}-\theta_{1}\right| \lambda_{m}^{2 m_{0}} \log \lambda_{m}\right. \\
& \left.+\left|d-d_{0}\right|\left|\theta_{20}-\theta_{2}\right| \lambda_{m}^{3 d_{0}} \log \lambda_{m}+\left|d-d_{0}\right|\left|\theta_{30}-\theta_{3}\right| \lambda_{m}^{1+2 d_{0}} \log \lambda_{m}\right)
\end{aligned}
$$

since $A_{j}(d)=\cos (d \pi / 2)\left(1+O\left(\lambda_{j}\right)\right)$ and similarly for $B_{j}\left(d_{0}\right)$. The required result for (30) comes by the difference between both. Finally the last three terms of the right hand side of the Lemma come similarly from (31).

The following lemma adapts Lemma 4 in Sun and Phillips (2003) to the non stationary case, allowing also for correlation between signal and noise.

Lemma 6 Let $0<d_{0}<3 / 4$ and $c_{k n}=c_{k}$ be a triangular array for which

$$
\max _{k}\left|c_{k}\right|=o(m), \quad \sum_{k=\left[1+m^{\alpha}\right]}^{m} c_{k}^{2} \sim \rho m, \quad \sum_{k=\left[1+m^{\alpha}\right]}^{m}\left|c_{k}\right|^{p}=O(m),
$$

for all $p \geq 1$ and $1>\alpha>\max \left\{1 / 2,1 / 4\left(1-d_{0}\right)\right\}$. Then

$$
\frac{1}{\sqrt{m}} \sum_{k=\left[1+m^{\alpha}\right]}^{m} c_{k} U_{z k} \xrightarrow{d} N\left(0, \frac{\pi^{2}}{6} \rho\right)
$$

Proof: Note first that such and $\alpha$ always exists because $d_{0}<3 / 4$. The main difference with respect to Lemma 4 in Sun and Phillips (2003) comes from the possibility of $d_{0} \geq 1 / 2$. In view of Lemma 1, in order to have the error terms in the covariance matrix of the normalized discrete Fourier transforms to be $o\left(m^{-1 / 2}\right)$ we need to consider only Fourier frequencies $\lambda_{k}$ for $m^{\alpha}<k \leq m$ (compare with the trimming in Velasco 1999). The result follows then as in Robinson (1995), Velasco (1999) or Sun and Phillips (2003).

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[^1]:    ${ }^{1}$ Proposition 4.1 of Hurvich et al. (2005) contains two typos. The element $(2,2)$ of $\Gamma^{*}$ should be divided

[^2]:    ${ }^{2}$ A limited Monte Carlo, not reported but available upon request, confirm that the performance of the MGSE improves significantly in those cases where truncation is not needed and the parameter space is correspondingly adjusted, e.g. positive correlation and the elements in $\theta$ restricted to be positive.

[^3]:    ${ }^{3}$ Note that there is no theoretical justification for the standard error of MGSEsin since Hurvich et al. (2005) only covers MGSE and MGSEcos. The approximation for the standard error of MGSEsin is used appealing to the comments in Remark 4.

