# A UNIFORM LAW FOR CONVERGENCE TO THE LOCAL TIMES OF LINEAR FRACTIONAL STABLE MOTIONS 

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#### Abstract

We provide a uniform law for the weak convergence of additive functionals of partial sum processes to the local times of linear fractional stable motions, in a setting sufficiently general for statistical applications. Our results are fundamental to the analysis of the global properties of nonparametric estimators of nonlinear statistical models that involve such processes as covariates.


1. Introduction. Let $x_{t}=\sum_{s=1}^{t} v_{s}$ be the partial sum of a scalar linear process $\left\{v_{t}\right\}$, for which the finite-dimensional distributions of $d_{n}^{-1} x_{\lfloor n r\rfloor}$ converge to those of $X(r)$. Under certain regularity conditions, we then have the finitedimensional convergence

$$
\begin{equation*}
\mathcal{L}_{n}^{f}\left(a, h_{n}\right):=\frac{d_{n}}{n h_{n}} \sum_{t=1}^{n} f\left(\frac{x_{t}-d_{n} a}{h_{n}}\right) \underset{\text { f.d.d. }}{\rightsquigarrow} \mathcal{L}(a) \int_{\mathbb{R}} f, \tag{1.1}
\end{equation*}
$$

where $a \in \mathbb{R}, f$ is Lebesgue integrable, $h_{n}=o\left(d_{n}\right)$ is a deterministic sequence, and $\mathcal{L}$ denotes the occupation density (or local time; see Remark 2.5 below) associated to $X$. Convergence results of this kind are particularly well documented in the case where $\left\{x_{t}\right\}$ is a random walk [see the monograph by Borodin and Ibragimov (1995)], and have more recently been extended to cover generating mechanisms that allow the increments of $\left\{x_{t}\right\}$ to exhibit significant temporal dependence [Jeganathan (2004), Wang and Phillips (2009a)].

These more general theorems concerning (1.1) have, in turn, played a fundamental role in the study of nonparametric estimation and testing in the setting of nonlinear cointegrating models. The simplest of these models takes the form

$$
\begin{equation*}
y_{t}=m_{0}\left(x_{t}\right)+u_{t}, \tag{1.2}
\end{equation*}
$$

where $\left\{x_{t}\right\}$ is as above, $\left\{u_{t}\right\}$ is a weakly dependent error process, and $m_{0}$ is an unknown function, assumed to possess a certain degree of smoothness (or be otherwise approximable). In a series of recent papers, (1.1) has facilitated the de-

[^0]velopment of a pointwise asymptotic distribution theory for kernel regression estimators of $m_{0}$ under very general conditions: see especially Wang and Phillips (2009a, 2009b, 2011, 2015), Kasparis and Phillips (2012) and Kasparis, Andreou and Phillips (2012). ${ }^{1}$

However, there are definite limits to the range of problems that can be successfully addressed with the aid of (1.1). In particular, since it concerns only the finite-dimensional convergence of $\mathcal{L}_{n}^{f}\left(a, h_{n}\right),(1.1)$ is suited only to studying the local behavior of a nonparametric estimator: that is, its behavior in the vicinity of a fixed spatial point. For the purpose of obtaining uniform rates of convergence for kernel regression estimators on "wide" domains-that is, on domains having a width of the same order as the range of $\left\{x_{t}\right\}_{t=1}^{n}$-it is manifestly inadequate. [See Duffy (2015), for a detailed account.] The situation is even worse with regard to sieve nonparametric estimation in this setting-which initially motivated the author's research on this problem-since in this case the development of even a pointwise asymptotic distribution theory requires a prior result on the uniform consistency of the estimator, over the entire domain on which estimation is to be performed.

The main purpose of this paper is thus to provide conditions under which the finite-dimensional convergence in (1.1) can be strengthened to the weak convergence

$$
\begin{equation*}
\mathcal{L}_{n}^{f}\left(a, h_{n}\right) \rightsquigarrow \mathcal{L}(a) \int_{\mathbb{R}} f, \tag{1.3}
\end{equation*}
$$

where $\mathcal{L}_{n}^{f}\left(a, h_{n}\right)$ is regarded as a process indexed by $(f, a) \in \mathscr{F} \times \mathbb{R}$, and $\left\{h_{n}\right\}$ may be random. Results of this kind are available in the existing literature, but only in the random walk case, which requires that the increments of $\left\{x_{t}\right\}$ be independent, and $X$ to be an $\alpha$-stable Lévy motion [see Borodin (1981, 1982); Perkins (1982); and Borodin and Ibragimov (1995), Chapter V]. In contrast, we allow the increments of $\left\{x_{t}\right\}$ to be serially correlated, such that the associated limiting process $X$ may be a linear fractional stable motion, which subsumes the $\alpha$-stable Lévy motion and fractional Brownian motion as special cases. Further, we permit the bandwidth sequence $\left\{h_{n}\right\}$ to be a random process, subject only to certain weak asymptotic growth conditions: this is of considerable utility in statistical applications, where the assumption that $\left\{h_{n}\right\}$ is a "given" deterministic sequence seems quite unrealistic. Crucial to the proof of (1.3) is a novel order estimate for $\mathcal{L}_{n}^{f}(a, 1)$ when $\int f=0$, which is of interest in its own right.

[^1]The remainder of this paper is organized as follows. Our assumptions on the data generating mechanism are described in Section 2. The main result (Theorem 3.1) is discussed in Section 3. An outline of the proof follows in Section 4, together with the statement of two key auxiliary results (Propositions 4.1 and 4.2). A preliminary application of our results to the kernel nonparametric estimation of $m_{0}$ in (1.2) is given in Section 5. The proof of Theorem 3.1 appears in Section 6, followed in Section 7 by proofs of Propositions 4.1 and 4.2. A proof related to the application appears in Section 8. The final two Sections 9 and 10 are of a more technical nature, detailing the proofs of two lemmas required in Section 7, and so may be skipped on a first reading.
1.1. Notation. For a complete listing of the notation used in this paper, see Section H of the Supplement. ${ }^{2}$ The stochastic order notations $o_{p}(\cdot)$ and $O_{p}(\cdot)$ have the usual definitions, as given, for example, in van der Vaart (1998), Section 2.2. For deterministic sequences $\left\{a_{n}\right\}$ and $\left\{b_{n}\right\}$, we write $a_{n} \sim b_{n}$ if $\lim _{n \rightarrow \infty} a_{n} / b_{n}=1$, and $a_{n} \asymp b_{n}$ if $\lim _{n \rightarrow \infty} a_{n} / b_{n} \in(-\infty, \infty) \backslash\{0\}$; for random sequences, $a_{n} \lesssim_{p} b_{n}$ denotes $a_{n}=O_{p}\left(b_{n}\right) . X_{n} \rightsquigarrow X$ denotes weak convergence in the sense of van der Vaart and Wellner (1996), and $X_{n} \rightsquigarrow$ f.d.d. $X$ the convergence of finite-dimensional distributions. For a metric space $(Q, d), \ell_{\infty}(Q)\left[\right.$ resp., $\left.\ell_{\text {ucc }}(Q)\right]$ denotes the space of uniformly bounded functions on $Q$, equipped with the topology of uniform convergence (resp., uniform convergence on compacta). For $p \geq 1, X$ a random variable, and $f: \mathbb{R} \rightarrow \mathbb{R},\|X\|_{p}:=\left(\mathbb{E}|X|^{p}\right)^{1 / p}$ and $\|f\|_{p}:=\left(\int_{\mathbb{R}}|f|^{p}\right)^{1 / p}$. BI denotes the space of bounded and Lebesgue integrable functions on $\mathbb{R} .\lfloor\cdot\rfloor$ and $\lceil\cdot\rceil$, respectively, denote the floor and ceiling functions. $C$ denotes a generic constant that may take different values even at different places in the same proof; $a \lesssim b$ denotes $a \leq C b$.
2. Model and assumptions. Our assumptions on the generating mechanism are similar to those of Jeganathan (2004)—who proves a finite-dimensional counterpart to our main theorem - and are comparable to those made on the regressor process in previous work on the estimation of nonlinear cointegrating regressions [see, e.g., Park and Phillips (2001), Wang and Phillips (2009b, 2012, 2015); and Kasparis and Phillips (2012)].

ASSUMPTION 1. (i) $\left\{\varepsilon_{t}\right\}$ is a scalar i.i.d. sequence. $\varepsilon_{0}$ lies in the domain of attraction of a strictly stable distribution with index $\alpha \in(0,2]$, and has characteristic function $\psi(\lambda):=\mathbb{E} \mathrm{e}^{\mathrm{i} \lambda \varepsilon_{0}}$ satisfying $\psi \in L^{p_{0}}$ for some $p_{0} \geq 1$.
(ii) $\left\{x_{t}\right\}$ is generated according to

$$
\begin{equation*}
x_{t}:=\sum_{s=1}^{t} v_{s}, \quad v_{t}:=\sum_{k=0}^{\infty} \phi_{k} \varepsilon_{t-k}, \tag{2.1}
\end{equation*}
$$

[^2]and either:
(a) $\alpha \in(1,2], \sum_{k=0}^{\infty}\left|\phi_{k}\right|<\infty$ and $\phi:=\sum_{k=0}^{\infty} \phi_{k} \neq 0$; or $\phi_{k} \sim k^{H-1-1 / \alpha} \pi_{k}$ for some $\left\{\pi_{k}\right\}_{k \geq 0}$ strictly positive and slowly varying at infinity, with
(b) $H>1 / \alpha$; or
(c) $H<1 / \alpha$ and $\sum_{k=0}^{\infty} \phi_{k}=0$.

In both cases (b) and (c), $H \in(0,1)$.
REMARK 2.1. Part (i) implies that there exists a slowly varying sequence $\left\{\varrho_{k}\right\}$ such that

$$
\begin{equation*}
\frac{1}{n^{1 / \alpha} \varrho_{n}} \sum_{t=1}^{\lfloor n r\rfloor} \varepsilon_{t} \underset{\text { f.d.d. }}{\rightsquigarrow} Z_{\alpha}(r), \tag{2.2}
\end{equation*}
$$

where $Z_{\alpha}$ denotes an $\alpha$-stable Lévy motion on $\mathbb{R}$, with $Z_{\alpha}(0)=0$. That is, the increments of $Z_{\alpha}$ are stationary, and for any $r_{1}<r_{2}$ the characteristic function of $Z_{\alpha}\left(r_{2}\right)-Z_{\alpha}\left(r_{1}\right)$ has the logarithm

$$
-\left(r_{2}-r_{1}\right) c|\lambda|^{\alpha}\left[1+\mathrm{i} \beta \operatorname{sgn}(\lambda) \tan \left(\frac{\pi \alpha}{2}\right)\right]
$$

where $\beta \in[-1,1]$ and $c>0$; following Jeganathan (2004), page 1773, we impose the further restriction that $\beta=0$ when $\alpha=1$. We shall also require that $\left\{\varrho_{k}\right\}$ be chosen such that $c=1$ here, which provides a convenient normalization for the scale of $Z_{\alpha}$.

REMARK 2.2. To permit the alternative forms of (ii) to be more concisely referenced, we shall refer to (a) as corresponding to the case where $H=1 / \alpha$; this designation may be justified by the manner in which the finite-dimensional limit of $d_{n}^{-1} x_{\lfloor n r\rfloor}$ depends on $(H, \alpha)$, as displayed in (2.6) below. The statement that $H<1 / \alpha$ will also be used as a shorthand for (c), that is, it will always be understood that $\sum_{k=0}^{\infty} \phi_{k}=0$ in this case.

We shall treat the parameters (including $H$ and $\alpha$ ) describing the data generating mechanism as "fixed" throughout, ignoring the dependence of any constants on these. Let $\left\{c_{k}\right\}$ denote a sequence with $c_{0}=1$ and

$$
c_{k}= \begin{cases}\phi, & \text { if } H=1 / \alpha  \tag{2.3}\\ |H-1 / \alpha|^{-1} k^{H-1 / \alpha} \pi_{k}, & \text { otherwise }\end{cases}
$$

By Karamata's theorem [Bingham, Goldie and Teugels (1987), Theorem 1.5.11], $\sum_{l=0}^{k} \phi_{k} \sim c_{k}$ as $k \rightarrow \infty$. Set

$$
\begin{equation*}
d_{k}:=k^{1 / \alpha} c_{k} \varrho_{k}, \quad e_{k}:=k d_{k}^{-1} \tag{2.4}
\end{equation*}
$$

and note that the sequences $\left\{c_{k}\right\},\left\{d_{k}\right\}$ and $\left\{e_{k}\right\}$ are regularly varying with indices $H-1 / \alpha, H$ and $1-H$, respectively. Theorems 5.1-5.3 in Kasahara and Maejima (1988) yield

Proposition 2.1. Under Assumption 1,

$$
\begin{equation*}
X_{n}(r):=\frac{1}{d_{n}} x_{\lfloor n r\rfloor} \underset{\text { f.d.d. }}{\rightsquigarrow} X(r), \quad r \in[0,1], \tag{2.5}
\end{equation*}
$$

where $X$ is the linear fractional stable motion (LFSM)

$$
\begin{align*}
X(r):= & \int_{0}^{r}(r-s)^{H-1 / \alpha} \mathrm{d} Z_{\alpha}(s)  \tag{2.6}\\
& +\int_{-\infty}^{0}\left[(r-s)^{H-1 / \alpha}-(-s)^{H-1 / \alpha}\right] \mathrm{d} Z_{\alpha}(s)
\end{align*}
$$

with the convention that $X=Z_{\alpha}$ when $H=1 / \alpha ; Z_{\alpha}$ is an $\alpha$-stable Lévy motion on $\mathbb{R}$, with $Z_{\alpha}(0)=0$.

REMARK 2.3. For a detailed discussion of the LFSM, see Samorodnitsky and Taqqu (1994). When $\alpha=2, Z_{\alpha}$ is a Brownian motion with variance 2; if additionally $H \neq 1 / \alpha, X$ is thus a fractional Brownian motion.

REMARK 2.4. Excepting such cases as the following:
(i) $\alpha \in(1,2], H>1 / \alpha$ [Astrauskas (1983), Theorem 2];
(ii) $\alpha=2, H=1 / \alpha$ and $\mathbb{E} \varepsilon_{0}^{2}<\infty$ [Hannan (1979)]; and
(iii) $\alpha=2, H<1 / \alpha$ and $\mathbb{E}\left|\varepsilon_{0}\right|^{q}<\infty$ for some $q>2$ [Davidson and de Jong (2000), Theorem 3.1];
it may not be possible to strengthen the convergence in (2.5) to weak convergence on $\ell_{\infty}[0,1]$. Weak convergence may hold, however, with respect to a weaker topology, and we shall be principally concerned with whether this topology is sufficiently strong that

$$
\begin{equation*}
\inf _{r \in[0,1]} X_{n}(r) \rightsquigarrow \inf _{r \in[0,1]} X(r), \quad \sup _{r \in[0,1]} X_{n}(r) \rightsquigarrow \sup _{r \in[0,1]} X(r), \tag{2.7}
\end{equation*}
$$

such as would follow from weak convergence in the Skorokhod $M_{1}$ topology [see Skorohod (1956), Section 2.2.10]. When $H=1 / \alpha$, sufficient conditions for this kind of convergence-which entail further restrictions on $\left\{\phi_{k}\right\}$ than are imposed here-are given in Avram and Taqqu (1992), Theorem 2 and Tyran-Kamińska (2010), Theorem 1 and Corollary 1. However, when $H<1 / \alpha$ and $\alpha \in(0,2)$, the sample paths of $X$ are unbounded, and thus (2.7) cannot possibly hold [see Samorodnitsky and Taqqu (1994), Example 10.2.5]. In any case, (2.7) is not necessary for the main results of this paper; it merely permits Theorem 3.1 below to take a slightly strengthened form.

REMARK 2.5. In consequence of Theorem 3(i) in Jeganathan (2004), the convergence in (2.5) occurs jointly with

$$
\mathcal{L}_{n}^{f}(a):=\frac{1}{e_{n}} \sum_{t=1}^{n} f\left(x_{t}-d_{n} a\right) \underset{\text { f.d.d. }}{\rightsquigarrow} \mathcal{L}(a) \int_{\mathbb{R}} f, \quad a \in \mathbb{R}
$$

for every $f \in \mathrm{BI}$. Here, $\{\mathcal{L}(a)\}_{a \in \mathbb{R}}$ denotes the occupation density (local time) of $X$, a process which, almost surely, has continuous paths and satisfies

$$
\begin{equation*}
\int_{\mathbb{R}} f(x) \mathcal{L}(x) \mathrm{d} x=\int_{0}^{1} f(X(r)) \mathrm{d} r \tag{2.8}
\end{equation*}
$$

for all Borel measurable and locally integrable $f$. (For the existence of $\mathcal{L}$, see Theorem 0 in Jeganathan (2004); the path continuity may be deduced from Theorem 3.1 below.)
3. A uniform law for the convergence to local time. Our main result concerns the convergence

$$
\begin{equation*}
\mathcal{L}_{n}^{f}\left(a, h_{n}\right):=\frac{1}{e_{n} h_{n}} \sum_{t=1}^{n} f\left(\frac{x_{t}-d_{n} a}{h_{n}}\right) \rightsquigarrow \mathcal{L}(a) \int_{\mathbb{R}} f \tag{3.1}
\end{equation*}
$$

where $\mathcal{L}_{n}^{f}\left(a, h_{n}\right)$ is regarded as a process indexed by $(f, a) \in \mathscr{F} \times \mathbb{R} .(\mathscr{F} \times \mathbb{R}$ is endowed with the product topology, $\mathscr{F}$ having the $L^{1}$ topology, and $\mathbb{R}$ the usual Euclidean topology.) $\left\{h_{n}\right\}$ is a measurable bandwidth sequence that may be functionally dependent on $\left\{x_{t}\right\}$, or indeed upon any other elements of the probability space; it is required only to satisfy:

ASSUMPTION 2. $h_{n} \in \mathscr{H}_{n}:=\left[\underline{h}_{n}, \bar{h}_{n}\right]$ with probability approaching 1 (w.p.a.1), where $\bar{h}_{n}=o\left(d_{n}\right)$ and $\underline{h}_{n}^{-1}=o\left(e_{n} \log ^{-2} n\right)$.

Define

$$
\begin{equation*}
\mathrm{BI}_{\beta}:=\left\{\left.f \in \mathrm{BI}\left|\int_{\mathbb{R}}\right| f(x)| | x\right|^{\beta} \mathrm{d} x<\infty\right\} \tag{3.2}
\end{equation*}
$$

and let $\mathrm{BIL}_{\beta}$ denote the subset of Lipschitz continuous functions in $\mathrm{BI}_{\beta}$. In order to state conditions on $\mathscr{F} \subset$ BI that are sufficient for (3.1) to hold, we first recall some definitions familiar from the theory of empirical processes. A function $F: \mathbb{R} \rightarrow \mathbb{R}_{+}$ is termed an envelope for $\mathscr{F}$, if $\sup _{f \in \mathscr{F}}|f(x)| \leq F(x)$ for every $x \in \mathbb{R}$. Given a pair of functions $l, u \in L^{1}$, define the bracket

$$
[l, u]:=\left\{f \in L^{1} \mid l(x) \leq f(x) \leq u(x), \forall x \in \mathbb{R}\right\}
$$

we say that $[l, u]$ is an $\varepsilon$-bracket if $\|u-l\|_{1}<\varepsilon$, and a continuous bracket if $l$ and $u$ are continuous. Let $N_{[]}^{*}(\varepsilon, \mathscr{F})$ denote the minimum number of continuous $\varepsilon$-brackets required to cover $\mathscr{F}$.

Assumption 3. (i) $\mathscr{F} \subset \mathrm{BI}$ has envelope $F \in \mathrm{BIL}_{\beta}$, for some $\beta>0$; and (ii) for each $\varepsilon>0, N_{[]}^{*}(\varepsilon, \mathscr{F})<\infty$.

We may now state our main result, the proof of which appears in Section 6.

Theorem 3.1. Suppose Assumptions 1-3 hold. Then:
(i) (3.1) holds in $\ell_{\mathrm{ucc}}(\mathscr{F} \times \mathbb{R})$;
and if additionally (2.7) holds, then
(ii) (3.1) holds in $\ell_{\infty}(\mathscr{F} \times \mathbb{R})$.

REMARK 3.1. The case where $h_{n}=1, \mathscr{F}=\{f\}$ and $\left\{x_{t}\right\}$ is a random walkwhich here entails $H=1 / \alpha$ and $\phi_{i}=0$ for all $i \geq 1$-has been studied extensively: see in particular Borodin (1981, 1982), Perkins (1982) and Borodin and Ibragimov (1995), Chapter V. In those works, it is proved (under these more restrictive assumptions on $\left\{x_{t}\right\}$ ) that

$$
\frac{1}{e_{n}} \sum_{t=1}^{\lfloor n r\rfloor} f\left(x_{t}-d_{n} a\right) \rightsquigarrow \mathcal{L}(a ; r) \int_{\mathbb{R}} f
$$

on $\ell_{\infty}(\mathbb{R} \times[0,1])$, where $\mathcal{L}(a ; r)$ denotes the local time of $X$ restricted to [0,r]. Theorem 3.1 could be very easily extended in this direction; we have refrained from doing so only to keep the paper to a reasonable length. The principal contribution of Theorem 3.1 is thus to extend this convergence in a direction more suitable for statistical applications, by allowing $\left\{v_{t}\right\}$ to be serially correlated and the bandwidth sequence $\left\{h_{n}\right\}$ to be data-dependent.

REMARK 3.2. After the manuscript of this paper had been completed, we obtained a copy of an unpublished manuscript by Liu, Chan and Wang (2014) in which, under rather different assumptions from those given here, a result similar to Theorem 3.1 is proved (for a fixed $f$ and a deterministic sequence $\left\{h_{n}\right\}$ ). Regarding the differences between our main result and their Theorem 2.1, we may note particularly their requirement that there exist a sequence of processes $\left\{X_{n}^{*}\right\}$ with $X_{n}^{*}={ }_{d} X$, and a $\delta>0$ such that

$$
\begin{equation*}
\sup _{r \in[0,1]}\left|X_{n}(r)-X_{n}^{*}(r)\right|=o_{\text {a.s. }}\left(n^{-\delta}\right) \tag{3.3}
\end{equation*}
$$

a condition which excludes a large portion of the processes considered in this paper, in view of Remark 2.4 above. [The availability of (3.3) permits these authors to prove their result by an argument radically different from that developed here.] On the other hand, our results do not subsume theirs, since these authors do not require $v_{t}$ to be a linear process.

Although Assumption 3 requires that $\mathscr{F}$ have a smooth envelope and smooth brackets, it is perfectly consistent with $\mathscr{F}$ containing discontinuous functions. Indeed, Assumption 3 is consistent with such cases as the following, as verified in Section A of the Supplement. [We expect that boundedness and $\int|f(x) \| x|^{\beta} \mathrm{d} x<$ $\infty$ could also be relaxed through the use of a suitable truncation argument, such as is employed in the proof of Theorem V.4.1 in Borodin and Ibragimov (1995).]

Example 3.1 (Single function). $\mathscr{F}=\{f\}$ where $f \in \mathrm{BI}_{\beta}$, and is majorised by another function $F \in \mathrm{BIL}_{\beta}$, in the sense that $|f(x)| \leq F(x)$ for all $x \in \mathbb{R}$. This obtains trivially if $f$ is itself in $\operatorname{BIL}_{\beta}$ (simply take $F(x):=|f(x)|$ ), but is also consistent with $f \in \mathrm{BI}_{\beta}$ having finitely many discontinuities (at the points $\left\{a_{k}\right\}_{k=1}^{K}$, where $\left.a_{k}<a_{k+1}\right)$, and being Lipschitz continuous on $\left(-\infty, a_{1}\right) \cup\left[a_{K}, \infty\right)$; all that is really necessary here is for $f$ to have one-sided Lipschitz approximants. Importantly, this includes the case where $f(x)=\mathbf{1}\{x \in I\}$ for any bounded interval $I$.

EXAmple 3.2 (Parametric family). $\mathscr{F}=\{g(x, \theta) \mid \theta \in \Theta\} \subset \mathrm{BIL}_{\beta}$, where $\Theta$ is compact, and there exists a $\tau \in(0,1]$ and a $\dot{g} \in \mathrm{BIL}_{\beta}$ such that

$$
\left|g(x, \theta)-g\left(x, \theta^{\prime}\right)\right| \leq \dot{g}(x)\left\|\theta-\theta^{\prime}\right\|^{\tau}
$$

for all $\theta, \theta^{\prime} \in \Theta$.
Example 3.3 (Smooth functions). $\mathscr{F}=\left\{f \in C^{\tau}(\mathbb{R})| | f \mid \leq F\right\}$, where $F \in$ $\mathrm{BIL}_{\beta}$ and

$$
C_{L}^{\tau}(\mathbb{R}):=\left\{f \in \mathrm{BI} \mid \exists C_{f}<L \text { s.t. }\left|f(x)-f\left(x^{\prime}\right)\right| \leq C_{f}\left|x-x^{\prime}\right|^{\tau} \forall x, x^{\prime} \in \mathbb{R}\right\}
$$

for some $\tau \in(0,1]$ and $L<\infty$.

## 4. Outline of proof and auxiliary results.

4.1. Outline of proof. The principal relationships between the results in this paper are summarized in Figure 1. The proof of Theorem 3.1, depicted in the top half of the figure, proceeds as follows. To reduce the difficulties arising by the randomness of $h=h_{n}$, we decompose

$$
\begin{equation*}
\mathcal{L}_{n}^{f}(a, h)=\mathcal{L}_{n}^{\varphi}(a) \int_{\mathbb{R}} f+\left[\mathcal{L}_{n}^{f}(a, h)-\mathcal{L}_{n}^{\varphi}(a) \int_{\mathbb{R}} f\right] \tag{4.1}
\end{equation*}
$$

where

$$
\begin{equation*}
\varphi(x):=(1-|x|) \mathbf{1}\{|x| \leq 1\} \tag{4.2}
\end{equation*}
$$

denotes the triangular kernel function, and $\mathcal{L}_{n}^{\varphi}(a):=\mathcal{L}_{n}^{\varphi}(a, 1)$. (This choice of $\varphi$ is made purely for convenience; any compactly supported Lipschitz function would serve our purposes equally well here.) It thus suffices to show that $\mathcal{L}_{n}^{\varphi} \rightsquigarrow \mathcal{L}$ in $\ell^{\infty}(\mathbb{R})$, and that the bracketed term on the right-hand side of (4.1) is uniformly negligible over $(f, h) \in \mathscr{F} \times \mathscr{H}_{n}$.

In view of Remark 2.5 above, the finite-dimensional distributions of $\mathcal{L}_{n}^{\varphi}$ converge to those of $\mathcal{L}$. The asymptotic tightness of $\mathcal{L}_{n}^{\varphi}$ will follow from the bound on the spatial increments

$$
\mathcal{L}_{n}^{\varphi}\left(a_{1}\right)-\mathcal{L}_{n}^{\varphi}\left(a_{2}\right)=\frac{1}{e_{n}} \sum_{t=1}^{n}\left[\varphi\left(x_{t}-d_{n} a_{1}\right)-\varphi\left(x_{t}-d_{n} a_{2}\right)\right]=: \frac{1}{e_{n}} \sum_{t=1}^{n} g_{1}\left(x_{t}\right)
$$



FIG. 1. Outline of proofs.
given in Proposition 4.1 below. The bracketed term on the right-hand side of (4.1) may be written as

$$
\begin{equation*}
\frac{1}{e_{n}} \sum_{t=1}^{n}\left[\frac{1}{h} f\left(\frac{x_{t}-d_{n} a}{h}\right)-\varphi\left(x_{t}-d_{n} a\right) \int_{\mathbb{R}} f\right]=: \frac{1}{e_{n}} \sum_{t=1}^{n} g_{2}\left(x_{t}\right) \tag{4.3}
\end{equation*}
$$

Control of (4.3) over progressively denser subsets of $\mathscr{F} \times \mathscr{H}_{n}$ is provided by Proposition 4.2 below; the conjunction of a bracketing argument and the continuity of the brackets suffices to extend this to the entirety of $\mathscr{F} \times \mathscr{H}_{n}$.

By construction, both $\int g_{1}=0$ and $\int g_{2}=0$. The proofs of Propositions 4.1 and 4.2 may therefore be approached in a unified way, through the analysis of sums of the form

$$
\begin{equation*}
\mathcal{S}_{n} g:=\sum_{t=1}^{n} g\left(x_{t}\right) \tag{4.4}
\end{equation*}
$$

where $g$ ranges over a class $\mathscr{G}$, all members of which have the property that $\int g=0$. Such functions are termed zero energy functions [Wang and Phillips (2011)]; we shall correspondingly term $\left\{\mathcal{S}_{n} g\right\}_{g \in \mathscr{G}}$ a zero energy process. Such processes are "centered" in the sense that $e_{n}^{-1 / 2} \mathcal{S}_{n} g$ converges weakly to a mixed Gaussian variate [Jeganathan (2008), Theorem 5]; whereas $e_{n}^{-1} \mathcal{S}_{n} g \rightsquigarrow \mathcal{L}(0) \int g$ if $\int g \neq 0$.

Equation (4.4) will be handled by decomposing $\mathcal{S}_{n} g$ as

$$
\mathcal{S}_{n} g=\sum_{k=0}^{n-1} \mathcal{M}_{n k} g+\mathcal{N}_{n} g
$$

where each $\mathcal{M}_{n k} g$ is a martingale; see (7.4) below. We provide order estimates for the sums of squares and conditional variances of the $\mathcal{M}_{n k} g$ 's (Lemma 7.3); by an application of either Burkholder's inequality, or a tail bound due to Bercu and Touati (2008), these translate into estimates for the $\mathcal{M}_{n k} g$ 's themselves. Propositions 4.1 and 4.2 then follow by standard arguments.
4.2. Key auxiliary results. To state these, we introduce the quantity

$$
\begin{equation*}
\|f\|_{[\beta]}:=\inf \left\{\left.c \in \mathbb{R}_{+}| | \hat{f}(\lambda)|\leq c| \lambda\right|^{\beta}, \forall \lambda \in \mathbb{R}\right\} \tag{4.5}
\end{equation*}
$$

for $f \in \mathrm{BI}, \beta \in(0,1]$, and $\hat{f}(\lambda):=\int \mathrm{e}^{\mathrm{i} \lambda x} f(x) \mathrm{d} x$. It is easily verified that $\|f\|_{[\beta]}$ is indeed a norm on the space $\mathrm{BI}_{[\beta]}:=\left\{f \in \mathrm{BI} \mid\|f\|_{[\beta]}<\infty\right\}$ (modulo equality almost everywhere). Some useful properties of $\|f\|_{[\beta]}$ are collected in Lemma 9.1 below; in particular, it is shown that $\mathrm{BI}_{[\beta]}$ contains all $f \in \mathrm{BI}_{\beta}$ for which $\int f=0$. Define

$$
\begin{equation*}
\bar{\beta}_{H}:=\frac{1-H}{2 H} \wedge 1 \tag{4.6}
\end{equation*}
$$

noting that $\bar{\beta}_{H} \in(0,1]$ for all $H \in(0,1)$, and let $\|\cdot\|_{\tau_{2 / 3}}$ denote the Orlicz norm associated to the convex and increasing function

$$
\tau_{2 / 3}(x):= \begin{cases}x(e-1), & \text { if } x \in[0,1]  \tag{4.7}\\ \mathrm{e}^{x^{2 / 3}}-1, & \text { if } x \in(1, \infty)\end{cases}
$$

[See van der Vaart and Wellner (1996), page 95 for the definition of an Orlicz norm.] A bound on the spatial increments of $\mathcal{L}_{n}^{\varphi}$ is given by:

Proposition 4.1. For every $\beta \in\left(0, \bar{\beta}_{H}\right)$, there exists $C_{\beta}<\infty$ such that

$$
\sup _{a_{1}, a_{2} \in \mathbb{R}}\left\|\mathcal{L}_{n}^{\varphi}\left(a_{1}\right)-\mathcal{L}_{n}^{\varphi}\left(a_{2}\right)\right\|_{\tau_{2 / 3}} \leq C_{\beta}\left|a_{1}-a_{2}\right|^{\beta}
$$

The next result shall be applied to prove that the recentered sums (4.3) are uniformly negligible. Since the order estimate given below is of interest in its own
right [see Duffy (2015), for an example of how it may be used to determine the uniform order of the first-order bias of a nonparametric regression estimator], we shall state it at a slightly higher level of generality than is needed for our purposes here. For $\mathscr{F} \subset \mathrm{BI}_{[\beta]}$, define

$$
\begin{equation*}
\delta_{n}(\beta, \mathscr{F}):=\|\mathscr{F}\|_{\infty}+e_{n}^{1 / 2}\left(\|\mathscr{F}\|_{1}+\|\mathscr{F}\|_{2}\right)+e_{n} d_{n}^{-\beta}\|\mathscr{F}\|_{[\beta]}, \tag{4.8}
\end{equation*}
$$

where $\|\mathscr{F}\|:=\sup _{f \in \mathscr{F}}\|f\|$.
Proposition 4.2. Suppose $\beta \in\left(0, \bar{\beta}_{H}\right)$ and $\mathscr{F}_{n} \subset \mathrm{BI}_{[\beta]}$ with $\# \mathscr{F}_{n} \lesssim n^{C}$. Then

$$
\begin{equation*}
\max _{f \in \mathscr{F}_{n}}\left|\mathcal{S}_{n} f\right| \lesssim_{p} \delta_{n}\left(\beta, \mathscr{F}_{n}\right) \log n \tag{4.9}
\end{equation*}
$$

If also $\left\|\mathscr{F}_{n}\right\|_{1} \lesssim 1,\left\|\mathscr{F}_{n}\right\|_{[\beta]}=o\left(d_{n}^{\beta}\right)$ and $\left\|\mathscr{F}_{n}\right\|_{\infty}=o\left(e_{n} \log ^{-2} n\right)$, then

$$
\max _{f \in \mathscr{F}_{n}}\left|\mathcal{S}_{n} f\right|=o_{p}\left(e_{n}\right)
$$

REMARK 4.1. As is clear from the proof, if $\beta \in\left[\bar{\beta}_{H}, 1\right]$ then (4.9) holds in a modified form, with $e_{n} d_{n}^{-\beta}\|\mathscr{F}\|_{[\beta]}$ in (4.8) being replaced by

$$
\left[\sum_{k=1}^{n-1} d_{k}^{-(1+\beta)}+e_{n}^{1 / 2} \sum_{k=1}^{n-1} k^{-1 / 2} d_{k}^{-(1+2 \beta) / 2}\right]\|\mathscr{F}\|_{[\beta]} .
$$

The proofs of Propositions 4.1 and 4.2 are given in Section 7 below.
5. A preliminary application to nonparametric regression. Suppose that we observe $\left\{\left(y_{t}, x_{t}\right)\right\}_{t=1}^{n}$ generated according to the nonlinear cointegrating regression model

$$
y_{t}=m_{0}\left(x_{t}\right)+u_{t},
$$

where $\left\{u_{t}\right\}$ is some weakly dependent disturbance process. As shown in Wang and Phillips (2009b), under suitable smoothness conditions the unknown function $m_{0}$ may be consistently estimated, at each fixed $x \in \mathbb{R}$, by the Nadaraya-Watson estimator

$$
\begin{align*}
\hat{m}(x) & =\frac{\sum_{t=1}^{n} K_{h_{n}}\left(x_{t}-x\right) y_{t}}{\sum_{t=1}^{n} K_{h_{n}}\left(x_{t}-x\right)} \\
& =m_{0}(x)+\frac{\sum_{t=1}^{n} K_{h_{n}}\left(x_{t}-x\right)\left[\left(m_{0}\left(x_{t}\right)-m_{0}(x)\right)+u_{t}\right]}{\sum_{t=1}^{n} K_{h_{n}}\left(x_{t}-x\right)} \tag{5.1}
\end{align*}
$$

where $K_{h}(u):=h^{-1} K\left(h^{-1} u\right)$, and $K \in \mathrm{BI}$ is a positive, mean-zero kernel with $\int_{\mathbb{R}} K=1$.

Now consider the problem of determining the rate at which $\hat{m}$ converges uniformly to $m_{0}$. As a first step, we would need to obtain the uniform rate of divergence of the denominator in (5.1); it is precisely this rate that the preceding results allow us to compute. By Theorem 3.1,

$$
\begin{equation*}
\frac{1}{e_{n}} \sum_{t=1}^{n} K_{h_{n}}\left(x_{t}-d_{n} a\right) \rightsquigarrow \mathcal{L}(a) \tag{5.2}
\end{equation*}
$$

in $\ell_{\infty}(\mathbb{R})$, provided $\mathbb{P}\left\{h_{n} \in \mathscr{H}_{n}\right\} \rightarrow 1$, and $K$ satisfies the requirements of Example 3.1 (which seems broad enough to cover any reasonable choice of $K$ ). Since $a \mapsto \mathcal{L}(a)$ is random-being dependent on the trajectory of the limiting process $X$-we now face the problem of identifying a sequence of sets on which the left side of (5.2) can be uniformly bounded away from zero. A natural candidate is

$$
A_{n}^{\varepsilon}:=\left\{x \in \mathbb{R} \mid \mathcal{L}_{n}\left(d_{n}^{-1} x\right) \geq \varepsilon\right\}
$$

where $\varepsilon>0 . \mathcal{L}_{n}$ is trivially bounded away from zero on this set, whence

$$
\begin{equation*}
\sup _{x \in A_{n}^{\varepsilon}}\left[\sum_{t=1}^{n} K_{h_{n}}\left(x_{t}-x\right)\right]^{-1} \lesssim_{p} e_{n}^{-1} \tag{5.3}
\end{equation*}
$$

More significantly, for any given $\delta>0$, we may choose $\varepsilon>0$ such that

$$
\begin{equation*}
\limsup _{n \rightarrow \infty} \mathbb{P}\left\{\frac{1}{n} \sum_{t=1}^{n} \mathbf{1}\left\{x_{t} \notin A_{n}^{\varepsilon}\right\} \geq \delta\right\} \leq \delta \tag{5.4}
\end{equation*}
$$

see Section 8 for the proof. That is, $\varepsilon>0$ may be chosen such that $A_{n}^{\varepsilon}$ contains as large a fraction of the observed trajectory $\left\{x_{t}\right\}_{t=1}^{n}$ as is desired, in the limit as $n \rightarrow \infty$. [Were we to allow $\varepsilon=\varepsilon_{n} \rightarrow 0$, we could permit $\delta=\delta_{n} \rightarrow 0$ here, but the order of (5.3) would necessarily be increased.]

Note that the sample-dependence of $A_{n}^{\varepsilon}$ is necessary for it fulfill two roles here, by being both "small" enough for (5.3) to hold, but also "large" enough to be consistent with (5.4). If $A_{n}^{\varepsilon}$ were replaced by a sequence of deterministic intervals (or sets, more generally), then the maintenance of (5.3) would necessarily come at the cost of violating (5.4). For example, when $x_{t}$ is a random walk with finite variance ( $\alpha=2$ ), the "widest" sequence of intervals $\left[-a_{n}, a_{n}\right]$ for which (5.3) holds is one for which $a_{n}=o\left(n^{1 / 2}\right)$ : but in consequence, the fraction of any trajectory $\left\{x_{t}\right\}_{t=1}^{n}$ falling within such an interval will converge to 0 , as $n \rightarrow \infty$ [see Remark 2.8 in Duffy (2015)].

In this respect, the availability of Theorem 3.1 allows us to improve upon the analysis provided in an earlier paper by Chan and Wang (2014)—and, in the random walk case, that of Gao et al. (2015)-who obtain uniform convergence rates for $\hat{m}_{n}$ on precisely such intervals [see Duffy (2015) for further details]. We expect that it would also play a similarly important role in the derivation of uniform convergence rates for series regression estimators in this setting, by ensuring the eigenvalues of the design matrix diverge at an appropriate rate, when attempting to estimate $m_{0}$ on a sequence of domains that contains most of the observed $\left\{x_{t}\right\}_{t=1}^{n}$.
6. Proof of Theorem 3.1. We shall prove only part (i) of Theorem 3.1 here; the relatively minor modifications required for the proof of part (ii) are detailed in Section B of the Supplement. Let $M<\infty$ be given; it suffices to prove that (3.1) holds in $\ell_{\infty}[-M, M]$. To simplify the exposition, we shall require that $h_{n} \in \mathscr{H}_{n}$ always; the proof in the general case (where this occurs w.p.a.1) requires no new ideas. The proof involves three steps:
(i) show that $\mathcal{L}_{n}^{\varphi}(a) \rightsquigarrow \mathcal{L}(a)$, using Proposition 4.1 ;
(ii) deduce $\mathcal{L}_{n}^{f}\left(a, h_{n}\right) \rightsquigarrow \mathcal{L}(a) \int_{\mathbb{R}} f$ for $f \in \mathrm{BIL}_{\beta}$, using a recentering, Proposition 4.2 and the Lipschitz continuity of $f$;
(iii) extend this to all $f \in \mathscr{F} \subset \mathrm{BI}$, where $\mathscr{F}$ satisfies Assumption 3, via a bracketing argument.
(i) Let $\varphi$ be the triangular kernel function, as defined in (4.2) above, and set $\beta_{0}:=\bar{\beta}_{H} / 2$. Recall that $\mathcal{L}_{n}^{f}(a):=\mathcal{L}_{n}^{f}(a, 1)$. By Proposition 4.1 and Theorem 2.2.4 in van der Vaart and Wellner (1996),

$$
\begin{aligned}
& \left\|\sup _{\left\{a, a^{\prime} \in M| | a-a^{\prime} \mid \leq \delta\right\}}\left|\mathcal{L}_{n}^{\varphi}\left(a^{\prime}\right)-\mathcal{L}_{n}^{\varphi}(a)\right|\right\|_{1} \\
& \quad \lesssim \int_{0}^{\delta} \log ^{3 / 2}\left(M \varepsilon^{-1 / \beta_{0}}\right) \mathrm{d} \varepsilon+\delta \log ^{3 / 2}\left(M \delta^{-2 / \beta_{0}}\right) \lesssim C_{M} \delta^{1 / 2}
\end{aligned}
$$

whence $\mathcal{L}_{n}^{\varphi}$ is tight in $\ell_{\infty}[-M, M]$. Thus, in view of Remark 2.5,

$$
\begin{equation*}
\mathcal{L}_{n}^{\varphi}(a) \rightsquigarrow \mathcal{L}(a) \tag{6.1}
\end{equation*}
$$

in $\ell_{\infty}[-M, M]$ [see van der Vaart and Wellner (1996), Example 2.2.12].
(ii) Now let $f \in \mathrm{BIL}_{\beta}$; we may without loss of generality take $f$ to be bounded by unity, with a Lipschitz constant of unity. For the subsequent argument, it will be more convenient to work with the inverse bandwidth $b:=h^{-1}$. Define

$$
\mathscr{B}_{n}:=\left\{h^{-1} \mid h \in \mathscr{H}_{n}\right\}=\left[\underline{b}_{n}, \bar{b}_{n}\right]:=\left[\bar{h}_{n}^{-1}, \underline{h}_{n}^{-1}\right]
$$

and let $f_{(a, b)}(x):=b f\left[b\left(x-d_{n} a\right)\right]$, for $(a, b) \in \mathbb{R} \times \mathbb{R}_{+}$. Take $C_{n}:=\left[-n^{\gamma}, n^{\gamma}\right] \times$ $\mathscr{B}_{n}$, let $\mathscr{C}_{n} \subset C_{n}$ be a lattice of mesh $n^{-\delta}$, and let $p_{n}(a, b)$ denote the projection of $(a, b) \in C_{n}$ onto a nearest neighbor in $\mathscr{C}_{n}$ (with some tie-breaking rule). The following is a straightforward consequence of the Lipschitz continuity of $f$ (see Section C of the Supplement for the proof).

Lemma 6.1. For every $\gamma \geq 1$, there exists $\delta>0$ such that

$$
\sup _{(a, b) \in C_{n}} \frac{1}{e_{n}} \sum_{t=1}^{n}\left|f_{(a, b)}\left(x_{t}\right)-f_{p_{n}(a, b)}\left(x_{t}\right)\right|=o_{p}(1)
$$

By taking $\gamma \geq 1$, we may ensure that $C_{n} \supset[-M, M] \times \mathscr{B}_{n}$, for all $n$ sufficiently large. Thus, for $\varphi_{(a)}:=\varphi_{(a, 1)}$ and $\mu_{f}:=\int_{\mathbb{R}} f$,

$$
\begin{align*}
& \sup _{(a, b) \in[-M, M] \times \mathscr{B}_{n}}\left|\mathcal{L}_{n}^{f}\left(a, b^{-1}\right)-\mu_{f} \mathcal{L}_{n}^{\varphi}(a)\right|  \tag{6.2}\\
& \leq \sup _{(a, b) \in C_{n}} \frac{1}{e_{n}}\left|\sum_{t=1}^{n}\left[f_{(a, b)}\left(x_{t}\right)-\mu_{f} \varphi_{(a)}\left(x_{t}\right)\right]\right| \\
& \leq \sup _{(a, b) \in \mathscr{C}_{n}} \frac{1}{e_{n}}\left|\sum_{t=1}^{n}\left[f_{(a, b)}\left(x_{t}\right)-\mu_{f} \varphi_{(a)}\left(x_{t}\right)\right]\right|+o_{p}(1) \\
& =\sup _{g \in \mathscr{G}_{n}} \frac{1}{e_{n}}\left|\sum_{t=1}^{n} g\left(x_{t}\right)\right|+o_{p}(1) \tag{6.3}
\end{align*}
$$

by Lemma 6.1, and we have defined $\mathscr{G}_{n}:=\left\{f_{(a, b)}-\mu_{f} \varphi_{(a)} \mid(a, b) \in \mathscr{C}_{n}\right\}$. It is readily verified that $\|g\|_{1}=1, \# \mathscr{G}_{n}=\# \mathscr{C}_{n} \lesssim n^{1+\gamma+2 \delta}$, and using Lemma 9.1(ii),

$$
\sup _{g \in \mathscr{G}_{n}}\|g\|_{[\beta]} \lesssim \underline{b}_{n}^{-\beta}=o\left(d_{n}^{\beta}\right), \quad \sup _{g \in \mathscr{G}_{n}}\|g\|_{\infty} \leq \bar{b}_{n} \lesssim e_{n} \log ^{-2} n .
$$

Thus, $\mathscr{G}_{n}$ satisfies the requirements of Proposition 4.2, whence (6.3) is $o_{p}(1)$. Hence, in view of (6.1),

$$
\begin{equation*}
\mathcal{L}_{n}^{f}\left(a, h_{n}\right) \rightsquigarrow \mathcal{L}(a) \int_{\mathbb{R}} f \tag{6.4}
\end{equation*}
$$

in $\ell_{\infty}[-M, M]$, for every $f \in \mathrm{BIL}_{\beta}$.
(iii) Finally, for $f \in \mathrm{BI}$ define the centered process

$$
v_{n}(f, a):=\mathcal{L}_{n}^{f}\left(a, h_{n}\right)-\mathcal{L}_{n}^{\varphi}(a) \int_{\mathbb{R}} f
$$

For a given $\varepsilon>0$, let $\left\{l_{k}, u_{k}\right\}_{k=1}^{K}$ denote a collection of continuous $L^{1}$ brackets that cover $\mathscr{F}$, with $\left\|u_{k}-l_{k}\right\|_{1}<\varepsilon$; the existence of these is guaranteed by Assumption 3. We first note (see Section C of the Supplement for the proof)

Lemma 6.2. Under Assumption 3, the brackets $\left\{l_{k}, u_{k}\right\}_{k=1}^{K}$ can be chosen so as to lie in $\mathrm{BIL}_{\beta}$.

For each $f \in \mathscr{F}$, there exists a $k \in\{1, \ldots, K\}$ such that $l_{k} \leq f \leq u_{k}, \int_{\mathbb{R}}\left(u_{k}-\right.$ f) $<\varepsilon$, and

$$
\begin{aligned}
v_{n}(f, a) & \leq \frac{1}{e_{n}} \sum_{t=1}^{n}\left[\frac{1}{h_{n}} u_{k}\left(\frac{x_{t}-d_{n} x}{h_{n}}\right)-\varphi\left(x_{t}\right) \int_{\mathbb{R}} f\right] \\
& \leq v_{n}\left(u_{k}, a\right)+\mathcal{L}_{n}^{\varphi}(a) \int_{\mathbb{R}}\left(u_{k}-f\right) .
\end{aligned}
$$

Taking suprema,

$$
\begin{aligned}
\sup _{(f, a) \in \mathscr{F} \times[-M, M]} v_{n}(f, a) & \leq \max _{1 \leq k \leq K} \sup _{a \in[-M, M]} v_{n}\left(u_{k}, a\right)+\varepsilon \sup _{a \in[-M, M]} \mathcal{L}_{n}^{\varphi}(a) \\
& =\varepsilon \sup _{a \in[-M, M]} \mathcal{L}_{n}^{\varphi}(a)+o_{p}(1)
\end{aligned}
$$

with the second equality following by (6.4), since we may take $u_{k} \in \mathrm{BIL}_{\beta}$ by Lemma 6.2. Applying a strictly analogous argument to the lower bracketing functions, $l_{k}$, we deduce that

$$
\begin{equation*}
\sup _{(f, a) \in \mathscr{F} \times[-M, M]}\left|v_{n}(f, a)\right| \leq \varepsilon \sup _{a \in[-M, M]}\left|\mathcal{L}_{n}^{\varphi}(a)\right|+o_{p}(1)=o_{p}(1) \tag{6.5}
\end{equation*}
$$

whence (3.1) holds in $\ell_{\infty}[-M, M]$, in view of (6.1).
7. Controlling the zero energy process. The proofs of Propositions 4.1 and 4.2 rely on a telescoping martingale decomposition similar to that used to prove maximal inequalities for mixingales [for a textbook exposition, see, e.g., Davidson (1994), Sections 16.2-16.3], which reduces $\mathcal{S}_{n} f$ to a sum of martingale components. In order to pass from control over each of these components to an order estimate for $\mathcal{S}_{n} f$ itself, we shall need the following results, the first of which is a straightforward consequence of Theorem 2.1 in Bercu and Touati (2008), and the second of which is well known. For a martingale $M:=\left\{M_{t}\right\}_{t=0}^{n}$ with associated filtration $\mathcal{G}:=\left\{\mathcal{G}_{t}\right\}_{t=0}^{n}$, define

$$
\begin{equation*}
[M]:=\sum_{t=1}^{n}\left(M_{t}-M_{t-1}\right)^{2}, \quad\langle M\rangle:=\sum_{t=1}^{n} \mathbb{E}\left[\left(M_{t}-M_{t-1}\right)^{2} \mid \mathcal{G}_{t-1}\right] . \tag{7.1}
\end{equation*}
$$

We say that $M$ is initialised at zero if $M_{0}=0$. Let $\|\cdot\|_{\tau_{1}}$ denote the Orlicz norm associated to $\tau_{1}(x):=\mathrm{e}^{x}-1$.

Lemma 7.1. Let $\left\{\Theta_{n}\right\}$ denote a sequence of index sets, and $\left\{K_{n}\right\}$ a real sequence such that $\# \Theta_{n}+K_{n} \lesssim n^{C}$. Suppose that for each $n \in \mathbb{N}, k \in\left\{1, \ldots, K_{n}\right\}$ and $\theta \in \Theta_{n}, M_{n k}(\theta)$ is a martingale, initialised at zero, for which

$$
\begin{equation*}
\omega_{n k}^{2}:=\max _{\theta \in \Theta_{n}}\left\{\left\|\left[M_{n k}(\theta)\right]\right\|_{\tau_{1}} \vee\left\|\left\{M_{n k}(\theta)\right\rangle\right\|_{\tau_{1}}\right\}<\infty . \tag{7.2}
\end{equation*}
$$

Then

$$
\max _{\theta \in \Theta_{n}}\left|\sum_{k=1}^{K_{n}} M_{n k}(\theta)\right| \lesssim_{p}\left(\sum_{k=1}^{K_{n}} \omega_{n k}\right) \log n .
$$

Lemma 7.2. Let $Z$ be a random variable. Then:
(i) $\|Z\|_{p} \lesssim p!^{1 / p} \sigma$ for all $p \in \mathbb{N}$, if and only if $\|Z\|_{\tau_{1}} \lesssim \sigma$;
(ii) $\|Z\|_{2 p} \lesssim(3 p)!^{1 / 2 p} \sigma$ for all $p \in \mathbb{N}$, if and only if $\|Z\|_{\tau_{2 / 3}} \lesssim \sigma$.

The proofs of Lemmas 7.1 and 7.2 appear in Section D of the Supplement.
7.1. The martingale decomposition. For a fixed $f \in \mathrm{BI}_{[\beta]}$, it follows from Lemma 9.3(ii) below and the reverse martingale convergence theorem [Hall and Heyde (1980), Theorem 2.6], that

$$
\left\|\mathbb{E}_{t} f\left(x_{t+k}\right)\right\|_{\infty} \lesssim d_{k}^{-(1+\beta)} \rightarrow 0, \quad \mathbb{E}_{t-k} f\left(x_{t}\right) \xrightarrow{p} \mathbb{E} f\left(x_{t}\right) \neq 0
$$

for each $t \geq 0$ as $k \rightarrow \infty$; here $\mathbb{E}_{t} f\left(x_{t+k}\right):=\mathbb{E}\left[f\left(x_{t+k}\right) \mid \mathcal{F}_{-\infty}^{t}\right]$, for $\mathcal{F}_{s}^{t}:=$ $\sigma\left(\left\{\varepsilon_{r}\right\}_{r=s}^{t}\right)$. Because $\left\{f\left(x_{t}\right)\right\}$ is asymptotically unpredictable only in the "forwards" direction, we truncate the "usual" decomposition at $t=0$, writing

$$
f\left(x_{t}\right)=\sum_{k=1}^{t}\left[\mathbb{E}_{t-k+1} f\left(x_{t}\right)-\mathbb{E}_{t-k} f\left(x_{t}\right)\right]+\mathbb{E}_{0} f\left(x_{t}\right)
$$

Performing this for each $1 \leq t \leq n$ gives

$$
\begin{aligned}
\sum_{t=1}^{n} f\left(x_{t}\right)= & \mathbb{E}_{0} f\left(x_{1}\right)+\left[f\left(x_{1}\right)-\mathbb{E}_{0} f\left(x_{1}\right)\right] \\
& +\mathbb{E}_{0} f\left(x_{2}\right)+\left[f\left(x_{2}\right)-\mathbb{E}_{1} f\left(x_{2}\right)\right]+\left[\mathbb{E}_{1} f\left(x_{2}\right)-\mathbb{E}_{0} f\left(x_{2}\right)\right]+\cdots \\
& +\mathbb{E}_{0} f\left(x_{n}\right)+\left[f\left(x_{n}\right)-\mathbb{E}_{n-1} f\left(x_{n}\right)\right]+\left[\mathbb{E}_{n-1} f\left(x_{n}\right)-\mathbb{E}_{n-2} f\left(x_{n}\right)\right] \\
& +\cdots+\left[\mathbb{E}_{1} f\left(x_{n}\right)-\mathbb{E}_{0} f\left(x_{n}\right)\right]
\end{aligned}
$$

Defining

$$
\begin{equation*}
\xi_{k t} f:=\mathbb{E}_{t} f\left(x_{t+k}\right)-\mathbb{E}_{t-1} f\left(x_{t+k}\right) \tag{7.3}
\end{equation*}
$$

and collecting terms appearing in the same "column" of the preceding display, we thus obtain

$$
\begin{align*}
\mathcal{S}_{n} f & =\sum_{t=1}^{n} f\left(x_{t}\right)=\sum_{t=1}^{n} \mathbb{E}_{0} f\left(x_{t}\right)+\sum_{k=0}^{n-1} \sum_{t=k+1}^{n}\left[\mathbb{E}_{t-k} f\left(x_{t}\right)-\mathbb{E}_{t-k-1} f\left(x_{t}\right)\right] \\
& =\sum_{t=1}^{n} \mathbb{E}_{0} f\left(x_{t}\right)+\sum_{k=0}^{n-1} \sum_{t=1}^{n-k}\left[\mathbb{E}_{t} f\left(x_{t+k}\right)-\mathbb{E}_{t-1} f\left(x_{t+k}\right)\right]  \tag{7.4}\\
& =\sum_{t=1}^{n} \mathbb{E}_{0} f\left(x_{t}\right)+\sum_{k=0}^{n-1} \sum_{t=1}^{n-k} \xi_{k t} f \\
& =\mathcal{N}_{n} f+\sum_{k=0}^{n-1} \mathcal{M}_{n k} f
\end{align*}
$$

where

$$
\mathcal{N}_{n} f:=\sum_{t=1}^{n} \mathbb{E}_{0} f\left(x_{t}\right), \quad \mathcal{M}_{n k} f:=\sum_{t=1}^{n-k} \xi_{k t} f .
$$

A bound for $\left\|\mathcal{N}_{n} f\right\|_{\infty}$ is provided by Lemma 9.3(ii) below. $\left\{\xi_{k t} f, \mathcal{F}_{-\infty}^{t}\right\}_{t=1}^{n-k}$, by construction, forms a martingale difference sequence for each $k$, and so control over each of the martingale "pieces" $\mathcal{M}_{n k} f$ will follow from control over

$$
\mathcal{U}_{n k} f:=\left[\mathcal{M}_{n k} f\right]=\sum_{t=1}^{n-k} \xi_{k t}^{2} f, \quad \mathcal{V}_{n k} f:=\left\langle\mathcal{M}_{n k} f\right\rangle=\sum_{t=1}^{n-k} \mathbb{E}_{t-1} \xi_{k t}^{2} f
$$

in combination with either Burkholder's inequality [Hall and Heyde (1980), Theorem 2.10] or Lemma 7.1 above, as appropriate.
7.2. Proofs of Propositions 4.1 and 4.2. Define

$$
\varsigma_{n}(\beta, f):=\|f\|_{\infty}+\|f\|_{1}+\|f\|_{[\beta]} \sum_{t=1}^{n} d_{t}^{-(1+\beta)}
$$

and

$$
\sigma_{n k}^{2}(\beta, f):= \begin{cases}\|f\|_{\infty}^{2}+\|f\|_{2}^{2} e_{n}, & \text { if } k \in\left\{0, \ldots, k_{0}\right\}, \\ e_{n}\left[k^{-1} d_{k}^{-(1+2 \beta)}\|f\|_{[\beta]}^{2}+\mathrm{e}^{-\gamma_{1} k}\|f\|_{1}^{2}\right], \\ & \text { if } k \in\left\{k_{0}+1, \ldots, n-1\right\} .\end{cases}
$$

The following provides the requisite control over the components of (7.4).
Lemma 7.3. For any $\beta \in[0,1]$,

$$
\begin{equation*}
\left\|\mathcal{N}_{n} f\right\|_{\infty} \lesssim \varsigma_{n}(\beta, f) \tag{7.5}
\end{equation*}
$$

and for all $0 \leq k \leq n-1$,

$$
\begin{equation*}
\left\|\mathcal{U}_{n k} f\right\|_{\tau_{1}} \vee\left\|\mathcal{V}_{n k} f\right\|_{\tau_{1}} \lesssim \sigma_{n k}^{2}(\beta, f) \tag{7.6}
\end{equation*}
$$

The proof of (7.6), in turn, relies upon the following.
Lemma 7.4. For every $k \in\{0, \ldots, n-1\}, t \in\{1, \ldots, n-k\}$ and $\beta \in(0,1]$

$$
\left\|\xi_{k t}^{2} f\right\|_{\infty}+\sum_{s=1}^{n-k-t}\left\|\mathbb{E}_{t} \xi_{k, t+s}^{2} f\right\|_{\infty} \lesssim \sigma_{n k}^{2}(\beta, f)
$$

The proofs of these results are deferred to Sections 9 and 10 . We shall also need the following, for which we recall the definition of $\delta_{n}(\beta, \mathscr{G})$ given in (4.8) above.

Lemma 7.5. If $\beta \in\left(0, \bar{\beta}_{H}\right)$ and $\mathscr{G} \subset \mathrm{BI}_{[\beta]}$, then there exists a $C_{\beta}<\infty$ such that

$$
\sup _{f \in \mathscr{G}} \varsigma_{n}(\beta, f)+\sum_{k=0}^{n-1} \sup _{f \in \mathscr{G}} \sigma_{n k}(\beta, f) \leq C_{\beta} \delta_{n}(\beta, \mathscr{G}) .
$$

The proof appears in Section D of the Supplement. We now turn to:
Proof of Proposition 4.1. Let $g \in \mathrm{BI}_{[\beta]}$. Burkholder's inequality, and Lemmas 7.2(i) and 7.3 and give

$$
\left\|\mathcal{M}_{n k} g\right\|_{2 p} \leq b_{2 p}^{1 / 2 p}\left\|\mathcal{U}_{n k} g\right\|_{p}^{1 / 2} \lesssim\left(b_{2 p} \cdot p!\right)^{1 / 2 p} \sigma_{n k}(\beta, g) \lesssim(3 p)!^{1 / 2 p} \sigma_{n k}(\beta, g)
$$

for every $p \in \mathbb{N}$, where $b_{2 p}$ depends on $p$ in the manner prescribed by Burkholder's inequality. Hence, $\left\|\mathcal{M}_{n k} g\right\|_{\tau_{2 / 3}} \lesssim \sigma_{n k}(\beta, g)$ by Lemma 7.2(ii). Then by (7.4), Lemma 7.3 and Lemma 7.5 (taking $\mathscr{G}=\{g\}$ )

$$
\begin{equation*}
\left\|\mathcal{S}_{n} g\right\|_{\tau_{2 / 3}} \leq\left\|\mathcal{N}_{n} g\right\|_{\infty}+\sum_{k=0}^{n-1}\left\|\mathcal{M}_{n k} g\right\|_{\tau_{2 / 3}} \leq C \delta_{n}(\beta, g) \tag{7.7}
\end{equation*}
$$

for some $C<\infty$ depending on $\beta$.
For $a_{1}, a_{2} \in \mathbb{R}$, set $\Delta:=\left|a_{1}-a_{2}\right|$ and define

$$
\varphi_{\left[a_{1}, a_{2}\right]}(x):=\varphi\left(x-d_{n} a_{1}\right)-\varphi\left(x-d_{n} a_{2}\right)
$$

Let $\beta \in\left(0, \bar{\beta}_{H}\right)$. Since $\varphi$ is bounded and Lipschitz,

$$
\left\|\varphi_{\left[a_{1}, a_{2}\right]}\right\|_{\infty} \leq\left(d_{n} \Delta\right) \wedge 1 \leq d_{n}^{\beta} \Delta^{\beta}
$$

and further, since $\varphi$ is bounded and compactly supported,

$$
\left\|\varphi_{\left[a_{1}, a_{2}\right]}\right\|_{p} \leq 2\left\|\varphi_{\left[a_{1}, a_{2}\right]}\right\|_{\infty}^{\beta}\left\|\varphi^{1-\beta}\right\|_{p} \lesssim d_{n}^{\beta} \Delta^{\beta}
$$

for $p \in\{1,2\}$. Finally, by Lemma 9.1(iii),

$$
\left\|\varphi_{\left[a_{1}, a_{2}\right]}\right\|_{[\beta]} \lesssim d_{n}^{\beta} \Delta^{\beta}
$$

Then by (7.7) and the definition of $\delta_{n}(\beta, \mathscr{F})$,

$$
\begin{aligned}
\left\|\mathcal{L}_{n}^{\varphi}\left(a_{1}\right)-\mathcal{L}_{n}^{\varphi}\left(a_{2}\right)\right\|_{\tau_{2 / 3}} & =e_{n}^{-1}\left\|\mathcal{S}_{n} \varphi_{\left[a_{1}, a_{2}\right]}\right\|_{\tau_{2 / 3}} \\
& \leq C\left(e_{n}^{-1 / 2} \cdot d_{n}^{\beta} \Delta^{\beta}+d_{n}^{-\beta} \cdot d_{n}^{\beta} \Delta^{\beta}\right) \\
& \lesssim C \Delta^{\beta},
\end{aligned}
$$

for some $C$ depending on $\beta$; here we have used the fact that since $\beta<\bar{\beta}_{H} \leq \frac{1-H}{2 H}$, $\left\{e_{n}^{-1 / 2} d_{n}^{\beta}\right\}$ is regularly varying with index $H\left(\beta-\frac{1-H}{2 H}\right)<0$.

Proof of Proposition 4.2. In view of Lemmas 7.3 and 7.5, we have

$$
\max _{f \in \mathscr{F}_{n}}\left|\mathcal{N}_{n} f\right| \leq \max _{f \in \mathscr{F}_{n}} \varsigma_{n}(\beta, f) \lesssim_{p} \delta_{n}\left(\beta, \mathscr{F}_{n}\right),
$$

and by an application of Lemma 7.1,

$$
\max _{f \in \mathscr{F}_{n}}\left|\sum_{k=0}^{n-1} \mathcal{M}_{n k} f\right| \lesssim_{p} \delta_{n}\left(\beta, \mathscr{F}_{n}\right) \log n .
$$

Thus, (4.9) follows from (7.4).

For the second part of the result, note that under the stated conditions on $\mathscr{F}_{n}$,

$$
\left\|\mathscr{F}_{n}\right\|_{2} \leq\left\|\mathscr{F}_{n}\right\|_{\infty}^{1 / 2}\left\|\mathscr{F}_{n}\right\|_{1}^{1 / 2}=o\left[e_{n}^{1 / 2} \log ^{-1} n\right]
$$

whence

$$
e_{n}^{-1} \delta_{n}\left(\beta, \mathscr{F}_{n}\right)=o_{p}\left(\log ^{-1} n\right)+d_{n}^{-\beta} o_{p}\left(d_{n}^{\beta}\right)=o_{p}(1)
$$

whereupon the result follows by (4.9).
8. Proof of (5.4). Let $\mu_{n}(a):=\frac{1}{n} \sum_{t=1}^{n} \mathbf{1}\left\{d_{n}^{-1} x_{t} \leq a\right\}$ and $\mu(a):=$ $\int_{-\infty}^{a} \mathcal{L}(x) \mathrm{d} x$. It is shown in Section E of the Supplement that

$$
\begin{equation*}
\mu_{n} \rightsquigarrow \mu \tag{8.1}
\end{equation*}
$$

in $\ell_{\infty}(\mathbb{R})$, jointly with the convergence in Theorem 3.1.
Let $T(x):=\mathbf{1}\{x<\varepsilon\}$ and $\mathcal{L}_{n}(a):=\mathcal{L}_{n}^{K}\left(a, h_{n}\right)$. We first note that

$$
\frac{1}{n} \sum_{t=1}^{n} \mathbf{1}\left\{x_{t} \notin A_{n}^{\varepsilon}\right\}=\frac{1}{n} \sum_{t=1}^{n} \mathbf{1}\left\{\mathcal{L}_{n}\left(d_{n}^{-1} x_{t}\right)<\varepsilon\right\}=\int_{\mathbb{R}} T\left(\mathcal{L}_{n}(a)\right) \mathrm{d} \mu_{n}(a)
$$

by definition of $A_{n}^{\varepsilon}$ and $\mu_{n}$. We shall now suppose that $\mathcal{L}_{n} \xrightarrow{\text { a.s. }} \mathcal{L}_{n}$ in $\ell_{\text {ucc }}(\mathbb{R})$, and $\mu_{n} \xrightarrow{\text { a.s. }} \mu$ in $\ell_{\infty}(\mathbb{R})$, as may be justified [in view of Theorem 3.1 and (8.1)] by Theorem 1.10.3 in van der Vaart and Wellner (1996); let $\Omega_{0} \subset \Omega$ denote a set, having $\mathbb{P} \Omega_{0}=1$, on which this convergence occurs. Define

$$
\bar{T}(x)= \begin{cases}1, & \text { if } x \leq \varepsilon \\ \varepsilon^{-1}(2 \varepsilon-x), & \text { if } x \in(\varepsilon, 2 \varepsilon) \\ 0, & \text { if } x \geq 2 \varepsilon\end{cases}
$$

Then, fixing an $\omega \in \Omega_{0}$,

$$
\int_{\mathbb{R}} T\left(\mathcal{L}_{n}^{\omega}(a)\right) \mathrm{d} \mu_{n}^{\omega}(a) \leq \int_{\mathbb{R}} \bar{T}\left(\mathcal{L}_{n}^{\omega}(a)\right) \mathrm{d} \mu_{n}^{\omega}(a)=\int_{\mathbb{R}} F_{n}(a) \mathrm{d} \mu_{n}^{\omega}(a)
$$

where $F_{n}(a):=\left(\bar{T} \circ \mathcal{L}_{n}^{\omega}\right)(a)$. Now let $[c, d]$ be chosen such that $\mu^{\omega}(d)-$ $\mu^{\omega}(c)<\varepsilon$. Since $\bar{T}$ is uniformly continuous, $F_{n}(a) \rightarrow F(a):=\left(\bar{T} \circ \mathcal{L}^{\omega}\right)(a)$ uniformly over $a \in[c, d]$, whence

$$
\begin{aligned}
\int_{\mathbb{R}} F_{n}(a) \mathrm{d} \mu_{n}^{\omega}(a) & \leq \int_{[c, d]^{c}} \mathrm{~d} \mu_{n}(a)+\int_{[c, d]} F(a) \mathrm{d} \mu_{n}^{\omega}(a)+\sup _{a \in[c, d]}\left|F_{n}(a)-F(a)\right| \\
& \rightarrow \varepsilon+\int_{[c, d]} F(a) \mathrm{d} \mu^{\omega}(a) \\
& \leq \varepsilon+\int_{\mathbb{R}} F(a) \mathrm{d} \mu^{\omega}(a),
\end{aligned}
$$

where the convergence follows by the Portmanteau theorem [van der Vaart and Wellner (1996), Theorem 1.3.4], since $F$ is continuous.

Thus,

$$
\begin{align*}
\limsup _{n \rightarrow \infty} \mathbb{P}\left\{\frac{1}{n} \sum_{t=1}^{n} \mathbf{1}\left\{x_{t} \notin A_{n}^{\varepsilon}\right\} \geq \delta\right\} & \leq \limsup _{n \rightarrow \infty} \mathbb{P}\left\{\int_{\mathbb{R}} \bar{T}\left(\mathcal{L}_{n}(a)\right) \mathrm{d} \mu_{n}(a) \geq \delta\right\} \\
& \leq \mathbb{P}\left\{\varepsilon+\int_{\mathbb{R}} \bar{T}(\mathcal{L}(a)) \mathrm{d} \mu(a) \geq \delta\right\}  \tag{8.2}\\
& \leq \mathbb{P}\left\{\varepsilon+\int_{\mathbb{R}} \mathbf{1}\{\mathcal{L}(a) \leq 2 \varepsilon\} \mathcal{L}(a) \mathrm{d} a \geq \delta\right\},
\end{align*}
$$

where noting that $\mathcal{L}$ is the density of $\mu$, the final inequality follows from

$$
\int_{\mathbb{R}} \bar{T}(\mathcal{L}(a)) \mathrm{d} \mu(a) \leq \int_{\mathbb{R}} 1\{\mathcal{L}(a) \leq 2 \varepsilon\} \mathrm{d} \mu(a)=\int_{\mathbb{R}} 1\{\mathcal{L}(a) \leq 2 \varepsilon\} \mathcal{L}(a) \mathrm{d} a .
$$

Finally,

$$
\varepsilon+\int_{\mathbb{R}} 1\{\mathcal{L}(a) \leq 2 \varepsilon\} \mathcal{L}(a) \mathrm{d} a \xrightarrow{\text { a.s. }} \int_{\mathbb{R}} 1\{\mathcal{L}(a)=0\} \mathcal{L}(a) \mathrm{d} a=0
$$

as $\varepsilon \rightarrow 0$, by dominated convergence, and so $\varepsilon>0$ may be chosen such that the right side of (8.2) is less than $\delta$.
9. Results preliminary to the proofs of Lemmas 7.3 and 7.4. Our arguments shall rely heavily on the use of the inverse Fourier transform to analyse objects of the form $\mathbb{E}_{t} f\left(x_{t+k}\right)$, similarly to Borodin and Ibragimov (1995), Jeganathan (2004, 2008) and Wang and Phillips (2009b, 2011). Provided that $f \in \mathrm{BI}$ and $Y$ has an integrable characteristic function $\psi_{Y}$, the "usual" inversion formula

$$
\begin{equation*}
\mathbb{E} f\left(y_{0}+Y\right)=\frac{1}{2 \pi} \int_{\mathbb{R}} \hat{f}(\lambda) \mathrm{e}^{-\mathrm{i} \lambda y_{0}} \mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda Y} \mathrm{~d} \lambda \tag{9.1}
\end{equation*}
$$

for $y_{0} \in \mathbb{R}$, is still valid, even when $\hat{f}(\lambda)=\int f(x) \mathrm{e}^{\mathrm{i} \lambda x} \mathrm{~d} x$ is not integrable; these conditions will always be met whenever the inversion formula is required below. The following provides some useful bounds for $\hat{f}$.

Lemma 9.1. For every $f \in \mathrm{BI}$ and $\beta \in(0,1]$,
(i) $|\hat{f}(\lambda)| \leq\left(|\lambda|^{\beta}\|f\|_{[\beta]}\right) \wedge\|f\|_{1}$;
(ii) if $\int f=0$, then

$$
\|f\|_{[\beta]} \leq 2^{1-\beta} \inf _{y \in \mathbb{R}} \int_{\mathbb{R}}|f(x-y)||x|^{\beta} \mathrm{d} x
$$

and so $\mathrm{BI}_{[\beta]} \supseteq\left\{f \in \mathrm{BI}_{\beta} \mid \int f=0\right\}$;
(iii) if $f(x):=g\left(x-a_{1}\right)-g\left(x-a_{2}\right)$ for some $a_{1}, a_{2} \in \mathbb{R}$, then

$$
\|f\|_{[\beta]} \leq 2^{1-\beta}\left|a_{1}-a_{2}\right|^{\beta}\|g\|_{1} .
$$

Let $\mathcal{F}_{s}^{t}:=\sigma\left(\left\{\varepsilon_{r}\right\}_{r=s}^{t}\right)$, noting that $\mathcal{F}_{s_{1}}^{s_{2}} \Perp \mathcal{F}_{s_{3}}^{s_{4}}$ for $s_{1} \leq s_{2}<s_{3} \leq s_{4}$. For $0<$ $s<t$, we shall have frequent recourse to the following decomposition:

$$
\begin{equation*}
x_{t}=\sum_{k=1}^{t} v_{t}=\sum_{k=1}^{t} \sum_{l=0}^{\infty} \phi_{l} \varepsilon_{k-l}=: x_{s-1, t}^{*}+\sum_{i=0}^{t-s} \varepsilon_{t-i} \sum_{j=0}^{i} \phi_{j}=: x_{s-1, t}^{*}+x_{s, t, t}^{\prime}, \tag{9.2}
\end{equation*}
$$

where $x_{s-1, t}^{*} \Perp x_{s, t, t}^{\prime}$ and $x_{s-1, t}^{*}$ is $\mathcal{F}_{-\infty}^{s-1}$-measurable. ${ }^{3}$ Defining $a_{i}:=\sum_{j=0}^{i} \phi_{j}$, we may further decompose $x_{s, t, t}^{\prime}$ as

$$
\begin{equation*}
x_{s, t, t}^{\prime}=\sum_{i=s}^{t} a_{t-i} \varepsilon_{i}=\sum_{i=s}^{r} a_{t-i} \varepsilon_{i}+\sum_{i=r+1}^{t} a_{t-i} \varepsilon_{i}=: x_{s, r, t}^{\prime}+x_{r+1, t, t}^{\prime} \tag{9.3}
\end{equation*}
$$

where $x_{s, r, t}^{\prime}$ is $\mathcal{F}_{s}^{r}$-measurable, and $x_{r+1, t, t}^{\prime}$ is $\mathcal{F}_{r+1}^{t}$-measurable. The following property of the coefficients $\left\{a_{i}\right\}$ is particularly important: there exist $0<\underline{a} \leq$ $\bar{a}<\infty$, and a $k_{0} \in \mathbb{N}$ such that

$$
\begin{equation*}
\underline{a} \leq \inf _{k_{0}+1 \leq k} \inf _{\lfloor k / 2\rfloor \leq l \leq k} c_{k}^{-1}\left|a_{l}\right| \leq \sup _{k_{0}+1 \leq k} \sup _{\lfloor k / 2\rfloor \leq l \leq k} c_{k}^{-1}\left|a_{l}\right| \leq \bar{a} \tag{9.4}
\end{equation*}
$$

This is an easy consequence of Karamata's theorem. Throughout the remainder of the paper, $k_{0}$ refers to the object of (9.4); it is also implicitly maintained $k_{0} \geq 8 p_{0}$ for $p_{0}$ as in Assumption 1(i).

Having decomposed $x_{t}$ into a sum of independent components, we shall proceed to control such objects as the right-hand side of (9.1) with the aid of Lemma 9.1 and the following, which provides bounds on integrals involving the characteristic functions of some of those components of $x_{t}$. Recall that Assumption 1(i) is equivalent to the statement that

$$
\begin{equation*}
\log \psi(\lambda)=-|\lambda|^{\alpha} G(\lambda)\left[1+\mathrm{i} \beta \operatorname{sgn}(\lambda) \tan \left(\frac{\pi \alpha}{2}\right)\right] \tag{9.5}
\end{equation*}
$$

for all $\lambda$ in a neighborhood of the origin, where $G$ is even and slowly varying at zero [see Ibragimov and Linnik (1971), Theorem 2.6.5]. Here, as throughout the remainder of this paper, a slowly varying (or regularly varying) function is understood to take only strictly positive values, and have the property that $G(\lambda)=$ $G(|\lambda|)$ for every $\lambda \in \mathbb{R}$.

Lemma 9.2. Let $p \in[0,5], q \in(0,2]$ and $z_{1}, z_{2} \in \mathbb{R}_{+}$. Then:
(i) there exists a $\gamma_{1}>0$ such that, for every $t \geq 0$ and $k \geq k_{0}+1$,

$$
\int_{\mathbb{R}}\left(z_{1}|\lambda|^{p} \wedge z_{2}\right)\left|\mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda x_{t+1, t+k, t+k}^{\prime}}\right| \mathrm{d} \lambda \lesssim z_{1} d_{k}^{-(1+p)}+z_{2} \mathrm{e}^{-\gamma_{1} k}
$$

[^3]and if $F(u) \asymp G^{p / \alpha}(u)$ as $u \rightarrow 0$,
\[

$$
\begin{aligned}
& \int_{\mathbb{R}}\left(z_{1}\left|a_{k}\right|^{p}|\lambda|^{p+q} F\left(a_{k} \lambda\right) \wedge z_{2}\right) \mid \mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda x_{t+1, t+k, t+k}^{\prime} \mid \mathrm{d} \lambda} \\
& \quad \lesssim z_{1} k^{-p / \alpha} d_{k}^{-(1+q)}+z_{2} \mathrm{e}^{-\gamma_{1} k}
\end{aligned}
$$
\]

(ii) for every $t \geq 1, k \geq k_{0}+1$ and $s \in\left\{k_{0}+1, \ldots, t\right\}$,

$$
\int_{\mathbb{R}}\left|\mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda x_{t-s+1, t-1, t+k}^{\prime}}\right| \mathrm{d} \lambda \lesssim \frac{c_{s}}{c_{k+s}} d_{s}^{-1}
$$

The preceding summarizes and refines some of the calculations presented on pages 15-21 of Jeganathan (2008). It further implies:

Lemma 9.3. Let $f \in \mathrm{BI}$. Then:
(i) for every $t \geq 0$ and $k \geq k_{0}+1$

$$
\mathbb{E}_{t}\left|f\left(x_{t+k}\right)\right| \lesssim d_{k}^{-1}\|f\|_{1}
$$

(ii) if in addition $f \in \mathrm{BI}_{[\beta]}$, then for every $t \geq 0$ and $k \geq k_{0}+1$,

$$
\left|\mathbb{E}_{t} f\left(x_{t+k}\right)\right| \lesssim \mathrm{e}^{-\gamma_{1} k}\|f\|_{1}+d_{k}^{-(1+\beta)}\|f\|_{[\beta]}
$$

For the next result, define

$$
\vartheta\left(z_{1}, z_{2}\right):=\mathbb{E}\left[\mathrm{e}^{-\mathrm{i} z_{1} \varepsilon_{0}}-\mathbb{E} \mathrm{e}^{-\mathrm{i} z_{1} \varepsilon_{0}}\right]\left[\mathrm{e}^{-\mathrm{i} z_{2} \varepsilon_{0}}-\mathbb{E} \mathrm{e}^{-\mathrm{i} z_{2} \varepsilon_{0}}\right] .
$$

LEMMA 9.4. Uniformly over $z_{1}, z_{2} \in \mathbb{R}$,

$$
\left|\vartheta\left(z_{1}, z_{2}\right)\right| \lesssim\left[\left|z_{1}\right|^{\alpha} \tilde{G}\left(z_{1}\right) \wedge 1\right]^{1 / 2}\left[\left|z_{2}\right|^{\alpha} \tilde{G}\left(z_{2}\right) \wedge 1\right]^{1 / 2}
$$

where $\tilde{G}(u) \asymp G(u)$ as $u \rightarrow 0$.
Proofs of (9.1), (9.4) and the preceding lemmas are given in Section F of the Supplement.

## 10. Proofs of Lemmas 7.3 and 7.4.

Proof of Lemma 7.3. By Lemma 9.3(ii),

$$
\begin{aligned}
\left|\mathcal{N}_{n} f\right| & \leq \sum_{t=1}^{k_{0}}\left|\mathbb{E}_{0} f\left(x_{t}\right)\right|+\sum_{t=k_{0}+1}^{n}\left|\mathbb{E}_{0} f\left(x_{t}\right)\right| \\
& \lesssim\|f\|_{\infty}+\sum_{t=k_{0}+1}^{n}\left[\mathrm{e}^{-\gamma_{1} t}\|f\|_{1}+d_{t}^{-(1+\beta)}\|f\|_{[\beta]}\right]
\end{aligned}
$$

whence (7.5). Regarding (7.6), it follows from repeated application of the law of iterated expectations that

$$
\begin{align*}
\mathbb{E}\left|\mathcal{V}_{n k} f\right|^{p} \leq & p!\cdot \sum_{t_{1}=1}^{n-k} \cdots \sum_{t_{p-1}=t_{p-2}}^{n-k} \mathbb{E}\left[\mathbb{E}_{t_{1}-1}\left(\xi_{k t_{1}}^{2} f\right) \cdots \mathbb{E}_{t_{p-1}-1}\left(\xi_{k t_{p-1}}^{2} f\right)\right]  \tag{10.1}\\
& \times\left(\left\|\xi_{k t_{p-1}}^{2} f\right\|_{\infty}+\sum_{s=1}^{n-k-t_{p-1}}\left\|\mathbb{E}_{t_{p-1}-1} \xi_{k, t_{p-1}+s}^{2} f\right\|_{\infty}\right)
\end{align*}
$$

more details of the calculations leading to (10.1) are given in Section G of the Supplement. By Lemma 7.4, the final term on the right is bounded by $C \sigma_{n k}^{2}(\beta, f)$. Proceeding inductively, we thus obtain

$$
\mathbb{E}\left|\mathcal{V}_{n k} f\right|^{p} \lesssim p!\cdot C^{p} \sigma_{n k}^{2 p}(\beta, f)
$$

whence the required bound follows by Lemma 7.2(i). An analogous argument yields the same bound for $\mathcal{U}_{n k} f$.

Proof of Lemma 7.4. We shall obtain the required bound for $\mathbb{E}_{t} \xi_{k, t+s}^{2} f$ by providing a bound for $\mathbb{E}_{t-s} \xi_{k t}^{2} f$ (for $s \in\{1, \ldots, t\}$ ) that depends only on $k$ and $s$ (and not $t$ ), separately considering the cases where:
(i) $k \in\left\{k_{0}+1, \ldots, n-t\right\}$; and
(ii) $k \in\left\{0, \ldots, k_{0}\right\}$.
(i) Recall the decomposition given in (9.2) and (9.3) above, applied here to reduce $x_{t+k}$ to a sum of independent pieces,

$$
\begin{aligned}
x_{t+k} & =x_{0, t+k}^{*}+x_{1, t-1, t+k}^{\prime}+x_{t, t, t+k}^{\prime}+x_{t+1, t+k, t+k}^{\prime} \\
& =x_{0, t+k}^{*}+x_{1, t-1, t+k}^{\prime}+a_{k} \varepsilon_{t}+x_{t+1, t+k, t+k}^{\prime}
\end{aligned}
$$

with the convention that $x_{1, t-1, t+k}^{\prime}=0$ if $t=1$, so that by Fourier inversion,

$$
\begin{align*}
\xi_{k t} f= & \mathbb{E}_{t} f\left(x_{t+k}\right)-\mathbb{E}_{t-1} f\left(x_{t+k}\right) \\
= & \frac{1}{2 \pi} \int_{\mathbb{R}} \hat{f}(\lambda) \mathrm{e}^{-\mathrm{i} \lambda x_{0, t+k}^{*}} \mathrm{e}^{-\mathrm{i} \lambda x_{1, t-1, t+k}^{\prime}}  \tag{10.2}\\
& \times\left[\mathrm{e}^{-\mathrm{i} \lambda a_{k} \varepsilon_{t}}-\mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda a_{k} \varepsilon_{t}}\right] \mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda x_{t+1, t+k, t+k}^{\prime}} \mathrm{d} \lambda .
\end{align*}
$$

Then

$$
\begin{align*}
\xi_{k t}^{2} f=\frac{1}{(2 \pi)^{2}} \iint_{\mathbb{R}^{2}} & \hat{f}\left(\lambda_{1}\right) \hat{f}\left(\lambda_{2}\right) \mathrm{e}^{-\mathrm{i}\left(\lambda_{1}+\lambda_{2}\right) x_{0, t+k}^{*}} \mathrm{e}^{-\mathrm{i}\left(\lambda_{1}+\lambda_{2}\right) x_{1, t-1, t+k}^{\prime}} \\
& \times\left[\mathrm{e}^{-\mathrm{i} \lambda_{1} a_{k} \varepsilon_{t}}-\mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda_{1} a_{k} \varepsilon_{t}}\right]\left[\mathrm{e}^{-\mathrm{i} \lambda_{2} a_{k} \varepsilon_{t}}-\mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda_{2} a_{k} \varepsilon_{t}}\right]  \tag{10.3}\\
& \times \mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda_{1} x_{t+1, t+k, t+k}^{\prime}} \mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda_{2} x_{t+1, t+k, t+k}^{\prime}} \mathrm{d} \lambda_{1} \mathrm{~d} \lambda_{2}
\end{align*}
$$

Now suppose $s \in\{k+1, \ldots, t\}$. Taking conditional expectations on both sides of (10.3) gives

$$
\begin{aligned}
\mathbb{E}_{t-s} \xi_{k t}^{2} f=\frac{1}{(2 \pi)^{2}} \iint_{\mathbb{R}^{2}} & \hat{f}\left(\lambda_{1}\right) \hat{f}\left(\lambda_{2}\right) \mathrm{e}^{-\mathrm{i}\left(\lambda_{1}+\lambda_{2}\right) x_{0, t+k}^{*}} \mathrm{e}^{-\mathrm{i}\left(\lambda_{1}+\lambda_{2}\right) x_{1, t-s, t+k}^{\prime}} \\
& \times \mathbb{E} \mathrm{e}^{-\mathrm{i}\left(\lambda_{1}+\lambda_{2}\right) x_{t-s+1, t-1, t+k}^{\prime}} \cdot \vartheta\left(\lambda_{1} a_{k}, \lambda_{2} a_{k}\right) \\
& \times \mathbb{E}^{-\mathrm{i} \lambda_{1} x_{t+1, t+k, t+k}^{\prime} \mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda_{2} x_{t+1, t+k, t+k}^{\prime}} \mathrm{d} \lambda_{1} \mathrm{~d} \lambda_{2}}
\end{aligned}
$$

where we have defined

$$
\vartheta\left(z_{1}, z_{2}\right):=\mathbb{E}\left[\mathrm{e}^{-\mathrm{i} z_{1} \varepsilon_{0}}-\mathbb{E} \mathrm{e}^{-\mathrm{i} z_{1} \varepsilon_{0}}\right]\left[\mathrm{e}^{-\mathrm{i} z_{1} \varepsilon_{0}}-\mathbb{E} \mathrm{e}^{-\mathrm{i} z_{2} \varepsilon_{0}}\right]
$$

for $z_{1}, z_{2} \in \mathbb{R}$, and made the further decomposition

$$
x_{1, t-1, t+k}^{\prime}=x_{1, t-s, t+k}^{\prime}+x_{t-s+1, t-1, t+k}^{\prime}
$$

with the convention that $x_{1, t-s, t+k}^{\prime}=0$ if $s=t$. Then, using (9.4), Lemma 9.4 and $|a b| \lesssim|a|^{2}+|b|^{2}$, we obtain

$$
\begin{aligned}
\mathbb{E}_{t-s} \xi_{k t}^{2} f \lesssim \iint_{\mathbb{R}^{2}} \mid & \hat{f}\left(\lambda_{1}\right) \hat{f}\left(\lambda_{2}\right) \mid \\
& \times\left[\left|a_{k} \lambda_{1}\right|^{\alpha} \tilde{G}\left(a_{k} \lambda_{1}\right) \wedge 1\right]^{1 / 2}\left[\left|a_{k} \lambda_{2}\right|^{\alpha} \tilde{G}\left(a_{k} \lambda_{2}\right) \wedge 1\right]^{1 / 2} \\
& \times\left|\mathbb{E} \mathrm{e}^{-\mathrm{i}\left(\lambda_{1}+\lambda_{2}\right) x_{t-s+1, t-1, t+k}^{\prime}}\right| \\
& \times\left|\mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda_{1} x_{t+1, t+k, t+k}^{\prime}}\right|\left|\mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda_{2} x_{t+1, t+k, t+k}^{\prime}}\right| \mathrm{d} \lambda_{1} \mathrm{~d} \lambda_{2} \\
\lesssim \int_{\mathbb{R}} \mid & \left.\hat{f}\left(\lambda_{1}\right)\right|^{2}\left(\left|a_{k}\right|^{\alpha}\left|\lambda_{1}\right|^{\alpha} \tilde{G}\left(a_{k} \lambda_{1}\right) \wedge 1\right)\left|\mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda_{1} x_{t+1, t+k, t+k}^{\prime}}\right|
\end{aligned}
$$

$$
\times \int_{\mathbb{R}}\left|\mathbb{E} \mathrm{e}^{-\mathrm{i}\left(\lambda_{1}+\lambda_{2}\right) x_{t-s+1, t-1, t+k}^{\prime}}\right| \mathrm{d} \lambda_{2} \mathrm{~d} \lambda_{1}
$$

where we have appealed to symmetry (in $\lambda_{1}$ and $\lambda_{2}$ ) to reduce the final bound to a single term. By a change of variables and Lemma 9.2(ii),

$$
\begin{align*}
\int_{\mathbb{R}} & \left|\mathbb{E} \mathrm{e}^{-\mathrm{i}\left(\lambda_{1}+\lambda_{2}\right) x_{t-s+1, t-1, t+k}^{\prime}}\right| \mathrm{d} \lambda_{2} \\
& =\int_{\mathbb{R}}\left|\mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda x_{t-s+1, t-1, t+k}^{\prime}}\right| \mathrm{d} \lambda \lesssim \frac{c_{s}}{c_{k+s}} d_{s}^{-1}, \tag{10.6}
\end{align*}
$$

while Lemma 9.1(i) and then Lemma 9.2(i) give

$$
\begin{align*}
& \int_{\mathbb{R}}|\hat{f}(\lambda)|^{2}\left(\left|a_{k}\right|^{\alpha}|\lambda|^{\alpha} \tilde{G}\left(a_{k} \lambda\right) \wedge 1\right) \mid \mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda x_{t+1, t+k, t+k}^{\prime} \mid \mathrm{d} \lambda} \\
& \quad \leq \int_{\mathbb{R}}\left[\left(\left|a_{k}\right|^{\alpha}|\lambda|^{\alpha+2 \beta} \tilde{G}\left(a_{k} \lambda\right)\|f\|_{[\beta]}^{2}\right) \wedge\|f\|_{1}^{2}\right] \mid \mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda x_{t+1, t+k, t+k}^{\prime} \mid \mathrm{d} \lambda}  \tag{10.7}\\
& \quad \lesssim k^{-1} d_{k}^{-(1+2 \beta)}\|f\|_{[\beta]}^{2}+\mathrm{e}^{-\gamma_{1} k}\|f\|_{1}^{2}
\end{align*}
$$

Together, (10.5)-(10.7) yield

$$
\begin{equation*}
\mathbb{E}_{t-s} \xi_{k t}^{2} f \lesssim \frac{c_{s}}{c_{k+s}} d_{s}^{-1}\left(k^{-1} d_{k}^{-(1+2 \beta)}\|f\|_{[\beta]}^{2}+\mathrm{e}^{-\gamma_{1} k}\|f\|_{1}^{2}\right) \tag{10.8}
\end{equation*}
$$

When $s \in\{1, \ldots, k\}$, (10.4) continues to hold, whence
(10.9)

$$
\begin{aligned}
& \mathbb{E}_{t-s} \xi_{k t}^{2} f \lesssim\left(\int_{\mathbb{R}}|\hat{f}(\lambda)|\left(|\lambda|^{\alpha / 2} \tilde{G}^{1 / 2}\left(a_{k} \lambda\right) \wedge 1\right) \mid \mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda x_{t+1, t+k, t+k}^{\prime} \mid \mathrm{d} \lambda}\right)^{2} \\
& \lesssim\left(\int_{\mathbb{R}}\left[\left(\left|a_{k}\right|^{\alpha / 2}|\lambda|^{(\alpha / 2+\beta)} \tilde{G}^{1 / 2}\left(a_{k} \lambda\right)\|f\|_{[\beta]}\right) \wedge\|f\|_{1}\right]\right. \\
&\left.\quad \times \mid \mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda x_{t+1, t+k, t+k}^{\prime} \mid \mathrm{d} \lambda}\right)^{2} \\
& \lesssim\left(k^{-1 / 2} d_{k}^{-(1+\beta)}\|f\|_{[\beta]}+\mathrm{e}^{-\gamma_{1} k}\|f\|_{1}\right)^{2} \\
& \lesssim d_{s}^{-1}\left(k^{-1} d_{k}^{-(1+2 \beta)}\|f\|_{[\beta]}^{2}+\mathrm{e}^{-\gamma_{1} k}\|f\|_{1}^{2}\right)
\end{aligned}
$$

by Lemmas 9.1(i) and 9.2(i); in obtaining the final result, we have used the fact that $s \leq k$ to replace a $d_{k}^{-1}$ by $d_{s}^{-1}$. Since $\left\{c_{k}\right\}$ is regularly varying and $k \geq k_{0}+1$, it follows from Potter's inequality [Bingham, Goldie and Teugels (1987), Theorem 1.5.6(iii)], that

$$
\sum_{s=1}^{k} d_{s}^{-1}+\sum_{s=k+1}^{n} \frac{c_{s}}{c_{k+s}} d_{s}^{-1} \lesssim \sum_{s=1}^{n} d_{s}^{-1} \lesssim n d_{n}^{-1}=e_{n}
$$

with the final bound following by Karamata's theorem. As noted above, since the bounds (10.8) and (10.9) do not depend on $t$, they apply also to $\mathbb{E}_{t} \xi_{k, t+s}^{2} f$. Hence, in view of the preceding,

$$
\begin{aligned}
\sum_{s=1}^{n-k-t} \mathbb{E}_{t} \xi_{k, t+s}^{2} f & \lesssim\left(k^{-1} d_{k}^{-(1+2 \beta)}\|f\|_{[\beta]}^{2}+\mathrm{e}^{-\gamma_{1} k}\|f\|_{1}^{2}\right)\left[\sum_{s=1}^{k} d_{s}^{-1}+\sum_{s=k+1}^{n-k-t} \frac{c_{s}}{c_{k+s}} d_{s}^{-1}\right] \\
& \lesssim e_{n}\left(k^{-1} d_{k}^{-(1+2 \beta)}\|f\|_{[\beta]}^{2}+\mathrm{e}^{-\gamma_{1} k}\|f\|_{1}^{2}\right)
\end{aligned}
$$

Turning now to $\left\|\xi_{k t}^{2} f\right\|_{\infty}$, note that (10.2) still holds, with the convention that $x_{1, t-1, t+k}=0$ if $t=1$. Thus, again by Lemmas 9.1(i) and 9.2(i),

$$
\begin{aligned}
\left\|\xi_{k t}^{2} f\right\|_{\infty} & \lesssim\left(\int_{\mathbb{R}}\left|\hat{f}(\lambda) \| \mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda x_{t+1, t+k, t+k}^{\prime}}\right| \mathrm{d} \lambda\right)^{2} \\
& \lesssim\left(\int_{\mathbb{R}}\left[|\lambda|^{\beta}\|f\|_{[\beta]} \wedge\|f\|_{1}\right]\left|\mathbb{E} \mathrm{e}^{-\mathrm{i} \lambda x_{t+1, t+k, t+k}^{\prime}}\right| \mathrm{d} \lambda\right)^{2} \\
& \lesssim\left(\|f\|_{[\beta]} d_{k}^{-(1+\beta)}+\|f\|_{1} \mathrm{e}^{-\gamma_{1} k}\right)^{2} \\
& \lesssim d_{k}^{-2(1+\beta)}\|f\|_{[\beta]}^{2}+\mathrm{e}^{-\gamma_{1} k}\|f\|_{1}^{2} \\
& \lesssim e_{n}\left(k^{-1} d_{k}^{-(1+2 \beta)}\|f\|_{[\beta]}^{2}+\mathrm{e}^{-\gamma_{1} k}\|f\|_{1}^{2}\right)
\end{aligned}
$$

where the final bound follows because $k \leq n$, and so $d_{k}^{-1} \lesssim k^{-1} n d_{n}^{-1}=k^{-1} e_{n}$.
(ii) When $s \in\left\{1, \ldots, k_{0}\right\}$, the crude bound $\mathbb{E}_{t-s} \xi_{k t}^{2} f \lesssim\|f\|_{\infty}^{2}$ suffices, since $k_{0}$ is fixed and finite. On the other hand, if $s \in\left\{k_{0}+1, \ldots, t\right\}$, we have by Jensen's inequality and Lemma 9.3(i) that

$$
\mathbb{E}_{t-s} \xi_{k t}^{2} f \leq \mathbb{E}_{t-s}\left(\mathbb{E}_{t} f\left(x_{t+k}\right)-\mathbb{E}_{t-1} f\left(x_{t+k}\right)\right)^{2} \lesssim \mathbb{E}_{t-s} f^{2}\left(x_{t+k}\right) \lesssim d_{s}^{-1}\|f\|_{2}^{2}
$$

Then, by Karamata's theorem,

$$
\begin{aligned}
\sum_{s=1}^{n-k-t} \mathbb{E}_{t} \xi_{k, t+s}^{2} f & \leq \sum_{s=1}^{k_{0}} \mathbb{E}_{t} \xi_{k, t+s}^{2} f+\sum_{s=k_{0}+1}^{n-k-t} \mathbb{E}_{t} \xi_{k, t+s}^{2} f \\
& \lesssim\|f\|_{\infty}^{2}+\|f\|_{2}^{2} \sum_{s=k_{0}+1}^{n-k-t} d_{s}^{-1} \\
& \lesssim\|f\|_{\infty}^{2}+\|f\|_{2}^{2} e_{n}
\end{aligned}
$$

Regarding $\left\|\xi_{k t}^{2} f\right\|_{\infty}$, the bound $\left\|\xi_{k t}^{2} f\right\|_{\infty} \lesssim\|f\|_{\infty}^{2}$ obtains trivially.
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## REFERENCES

Astrauskas, A. (1983). Limit theorems for sums of linearly generated random variables. Lith. Math. J. 23 127-134. MR0706002
AVram, F. and TAQQU, M. S. (1992). Weak convergence of sums of moving averages in the $\alpha$ stable domain of attraction. Ann. Probab. 20 483-503. MR1143432
Bercu, B. and Touati, A. (2008). Exponential inequalities for self-normalized martingales with applications. Ann. Appl. Probab. 18 1848-1869. MR2462551
Bingham, N. H., Goldie, C. M. and Teugels, J. L. (1987). Regular Variation. Cambridge Univ. Press, Cambridge. MR1015093
Borodin, A. N. (1981). The asymptotic behavior of local times of recurrent random walks with finite variance. Theory Probab. Appl. 26 758-772. MR0636771
Borodin, A. N. (1982). On the asymptotic behavior of local times of recurrent random walks with infinite variance. Theory Probab. Appl. 29 318-333. MR0749918
Borodin, A. N. and Ibragimov, I. A. (1995). Limit theorems for functionals of random walks. Proc. Steklov Inst. Math. 195 1-259. MR1368394
ChAN, N. and WANG, Q. (2014). Uniform convergence for nonparametric estimators with nonstationary data. Econometric Theory 30 1110-1133. MR3252034
DaVidson, J. (1994). Stochastic Limit Theory: An Introduction for Econometricians. Oxford Univ. Press, New York. MR1430804
Davidson, J. and de Jong, R. M. (2000). The functional central limit theorem and weak convergence to stochastic integrals. II. Fractionally integrated processes. Econometric Theory 16 643-666. MR1802836

DUFFY, J. A. (2015). Uniform convergence rates over maximal domains in structural nonparametric cointegrating regression. Unpublished manuscript, Univ. Oxford.
Gao, J., Kanaya, S., Li, D. and TjøStheim, D. (2015). Uniform consistency for nonparametric estimators in null recurrent time series. Econometric Theory 31 911-952. MR3396235
Hall, P. and Heyde, C. C. (1980). Martingale Limit Theory and Its Application (Probability and Mathematical Statistics). Academic Press, New York. MR0624435
Hannan, E. J. (1979). The central limit theorem for time series regression. Stochastic Process. Appl. 9 281-289. MR0562049
Ibragimov, I. A. and Linnik, Yu. V. (1971). Independent and Stationary Sequences of Random Variables. Wolters-Noordhoff Publishing, Groningen. MR0322926
Jeganathan, P. (2004). Convergence of functionals of sums of r.v.s to local times of fractional stable motions. Ann. Probab. 32 1771-1795. MR2073177
Jeganathan, P. (2008). Limit theorems for functionals of sums that converge to fractional Brownian and stable motions. Cowles Foundation Discussion Paper No. 1649, Yale Univ., New Haven, CT.
Karlsen, H. A., Myklebust, T. and Tuøstheim, D. (2007). Nonparametric estimation in a nonlinear cointegration type model. Ann. Statist. 35 252-299. MR2332276
Karlsen, H. A. and TıøSTHEIM, D. (2001). Nonparametric estimation in null recurrent time series. Ann. Statist. 29 372-416. MR1863963
Kasahara, Y. and Maejima, M. (1988). Weighted sums of i.i.d. random variables attracted to integrals of stable processes. Probab. Theory Related Fields 78 75-96. MR0940869
Kasparis, I., Andreou, E. and Phillips, P. C. B. (2012). Nonparametric predictive regression. Cowles Foundation Discussion Paper No. 1878, Yale Univ., New Haven, CT.
Kasparis, I. and Phillips, P. C. B. (2012). Dynamic misspecification in nonparametric cointegrating regression. J. Econometrics 168 270-284. MR2923768
LiU, W., Chan, N. and WANG, Q. (2014). Uniform approximation to local time with applications in non-linear cointegrating regression. Unpublished manuscript, Univ. Sydney.
Park, J. Y. and Phillips, P. C. B. (2001). Nonlinear regressions with integrated time series. Econometrica 69 117-161. MR1806536
Perkins, E. (1982). Weak invariance principles for local time. Probab. Theory Related Fields $\mathbf{6 0}$ 437-451. MR0665738
Samorodnitsky, G. and TaqQu, M. S. (1994). Stable Non-Gaussian Random Processes: Stochastic Models With Infinite Variance. Chapman \& Hall, New York. MR1280932
SKOROHOD, A. V. (1956). Limit theorems for stochastic processes. Theory Probab. Appl. 1 261290. MR0084897

Tyran-Kamińska, M. (2010). Functional limit theorems for linear processes in the domain of attraction of stable laws. Statist. Probab. Lett. 80 975-981. MR2638967
van der Vaart, A. W. (1998). Asymptotic Statistics. Cambridge Univ. Press, Cambridge. MR1652247
VAn der Vaart, A. W. and Wellner, J. A. (1996). Weak Convergence and Empirical Processes: With Applications to Statistics. Springer, New York. MR1385671
Wang, Q. and Phillips, P. C. B. (2009a). Asymptotic theory for local time density estimation and nonparametric cointegrating regression. Econometric Theory 25 710-738. MR2507529
WANG, Q. and Phillips, P. C. B. (2009b). Structural nonparametric cointegrating regression. Econometrica 77 1901-1948. MR2573873
Wang, Q. and Phillips, P. C. B. (2011). Asymptotic theory for zero energy functionals with nonparametric regression applications. Econometric Theory 27 235-259. MR2782038
Wang, Q. and Phillips, P. C. B. (2012). A specification test for nonlinear nonstationary models. Ann. Statist. 40 727-758. MR2933664

WANG, Q. and Phillips, P. C. B. (2015). Nonparametric cointegrating regression with endogeneity and long memory. Econometric Theory. To appear. DOI:10.1017/S0266466614000917.

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[^1]:    ${ }^{1}$ If $\left\{x_{t}\right\}$ is Markov, then this distribution theory may be developed by quite different arguments, without the use of (1.1), see Karlsen and Tjøstheim (2001) and Karlsen, Myklebust and Tjøstheim (2007), which have spawned a large literature. While we consider this approach to the problem to be equally important, our results touch upon it only a little, since we work with a class of regressor processes that are typically (excepting the random walk case) non-Markov.

[^2]:    ${ }^{2}$ The Supplement is available as an addendum to arXiv:1501.05467.

[^3]:    ${ }^{3} x_{s-1, t}^{*}$ is weighted sum of $\left\{\varepsilon_{t}\right\}_{t=-\infty}^{s-1}$ : since these weights are not important for our purposes, we have refrained from giving an explicit formula for these here.

