SPLIT INVARIANCE PRINCIPLES FOR STATIONARY PROCESSES

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The results of Komlós, Major and Tusnády give optimal Wiener approximation of partial sums of i.i.d. random variables and provide an extremely powerful tool in probability and statistical inference. Recently Wu [Ann. Probab. 35 (2007) 2294–2320] obtained Wiener approximation of a class of dependent stationary processes with finite pth moments, 2 ≤ p ≤ 4, with error term o(n^{1/p}(log n)^γ), γ > 0, and Liu and Lin [Stochastic Process. Appl. 119 (2009) 249–280] removed the logarithmic factor, reaching the Komlós–Major–Tusnády bound o(n^{1/p}). No similar results exist for p > 4, and in fact, no existing method for dependent approximation yields an a.s. rate better than o(n^{1/4}). In this paper we show that allowing a second Wiener component in the approximation, we can get rates near to o(n^{1/p}) for arbitrary p > 2. This extends the scope of applications of the results essentially, as we illustrate it by proving new limit theorems for increments of stochastic processes and statistical tests for short term (epidemic) changes in stationary processes. Our method works under a general weak dependence condition covering wide classes of linear and nonlinear time series models and classical dynamical systems.

1. Introduction. Let X, X₁, X₂, ... be i.i.d. random variables with mean 0 and variance 1, and let Sn = ∑k≤nXk. Komlós, Major and Tusnády [25, 26] showed that if E(e^{t|X|}) < ∞ for some t > 0 then, after suitably enlarging the probability space, there exists a Wiener process \{W(t), t ≥ 0\} such that

\[ \sum_{k≤n} X_k = W(n) + O((log n)^γ) \quad a.s. \]
Also, if \( E|X|^p < \infty \) for some \( p > 2 \), they proved the approximation
\[
S_n = W(n) + o(n^{1/p}) \quad \text{a.s.}
\]  
(2)

The remainder terms in (1) and (2) are optimal. In the case when only \( EX = 0, EX^2 = 1 \) is assumed, Strassen [46] obtained
\[
S_n = W(n) + o((n \log \log n)^{1/2}) \quad \text{a.s.}
\]  
(3)

Without additional moment assumptions the rate in (3) is also optimal (see Major [29]). Relation (3) is a useful invariance principle for the law of the iterated logarithm; on the other hand, it does not imply the CLT for \( \{X_n\} \).

This difficulty was removed by Major [30] who showed that under \( EX = 0, EX^2 = 1 \) there exists a Wiener process \( W \) and a numerical sequence \( \tau_n \sim n \) such that
\[
S_n = W(\tau_n) + o(n^{1/2}) \quad \text{a.s.}
\]  
(4)

Thus allowing a slight perturbation of the approximating Wiener process one can reach the remainder term \( o(n^{1/p}) \) also for \( p = 2 \), making the result applicable for a wide class of CLT-type results. The case of strong approximation under the moment condition \( EX^2h(|X|) < \infty \) where \( h(x) = o(x^\varepsilon) \), \( x \to \infty \), for any \( \varepsilon > 0 \), has been cleared up completely by Einmahl [18].

The previous results, which settle the strong approximation problem for i.i.d. random variables with finite variances, provide powerful tools in probability and statistical inference (see, e.g., Shorack and Wellner [45]). Starting with Strassen [47], a wide literature has dealt with extensions of the above results for weakly dependent sequences, but the existing results are much weaker than in the i.i.d. case. Recently, however, Wu [50] showed that for a large class of weakly dependent stationary sequences \( \{X_n\} \) satisfying \( E|X_1|^p < \infty, 2 < p \leq 4 \), we have the approximation
\[
S_n = W(n) + o(n^{1/p}(\log n)^{\gamma}) \quad \text{a.s.}
\]

for some \( \gamma > 0 \), and Liu and Lin [28] removed the logarithmic factor in the error term, reaching the optimal Komlós–Major–Tusnády bound. The proofs do not work for \( p > 4 \), and in fact, no existing method for dependent approximation yields an a.s. rate better than \( o(n^{1/4}) \). On the other hand, many important limit theorems in probability and statistics involve norming sequences smaller than \( n^{1/4} \), making such results inaccessible by invariance methods. The purpose of the present paper is to fill this gap and provide a new type of approximation theorem reaching nearly the Komlós–Major–Tusnády rate for any \( p > 2 \).

As noted above, reaching the error term \( o(n^{1/2}) \) for i.i.d. sequences with finite variance requires a perturbation of the approximating Wiener process \( W \). In the case of dependent processes we will also need a similar perturbation, and, more essentially, we will include a second Wiener process in
the approximation, whose scaling factor is smaller than that of $W$, and thus it will not affect the asymptotic behavior of the main term. Specifically, for a large class of weakly dependent stationary processes $\{Y_k\}$ with finite $p$th moments, $2 < p < \infty$, we will prove the approximation

$$
\sum_{k=1}^{n} Y_k = W_1(s_n^2) + W_2(t_n^2) + O(n^{(1+n)/p}) \quad \text{a.s.,}
$$

where $\{W_1(t), t \geq 0\}$ and $\{W_2(t), t \geq 0\}$ are standard Wiener processes, and $s_n, t_n$ are numerical sequences with

$$
s_n^2 \sim \sigma^2 n, \quad t_n^2 \sim cn^\gamma
$$

for some $0 < \gamma < 1$, $\sigma^2 > 0$, $c > 0$. The new element in (5) is the term $W_2(t_n^2)$ which, by its smaller scaling, does not disturb the asymptotic properties of $W_1(s_n^2)$. Note that the processes $W_1, W_2$ are not independent, but this will not present any difficulties in applications. (See also Proposition 1 in the next section.) The number $\eta$ depends on the weak dependence rate of $\{Y_k\}$ (introduced below), and can be made arbitrarily small under suitable rate conditions.

For $p > 0$ and a random variable $Y$, let $\|Y\|_p = (E|Y|^p)^{1/p}$. If $A$ and $B$ are subsets of $\mathbb{Z}$, we let $d(A, B) = \inf\{|a - b| : a \in A, b \in B\}$.

**Definition 1.** Let $\{Y_k, k \in \mathbb{Z}\}$ be a stochastic process, let $p \geq 1$ and let $\delta(m) \to 0$. We say that $\{Y_k, k \in \mathbb{Z}\}$ is weakly $\mathcal{M}$-dependent in $L^p$ with rate function $\delta(\cdot)$ if:

(A) For any $k \in \mathbb{Z}, m \in \mathbb{N}$ one can find a random variable $Y_k^{(m)}$ with finite $p$th moment such that

$$
\|Y_k - Y_k^{(m)}\|_p \leq \delta(m).
$$

(B) For any disjoint intervals $I_1, \ldots, I_r$ $(r \in \mathbb{N})$ of integers and any positive integers $m_1, \ldots, m_r$, the vectors $\{Y_j^{(m_1)}, j \in I_1\}, \ldots, \{Y_j^{(m_r)}, j \in I_r\}$ are independent provided $d(I_k, I_l) > \max\{m_k, m_l\}$ for $1 \leq k < l \leq r$.

We remark that our dependence condition is naturally preserved under smooth transformations. For example, if $\{Y_k\}$ is weakly $\mathcal{M}$-dependent in $L^p$ with rate $\delta(\cdot)$, and $h$ is a Lipschitz $\alpha$ function $(0 < \alpha \leq 1)$ with Lipschitz constant $K$, then by the monotonicity of $\|Y_k - Y_k^{(m)}\|_p$ in $p$ we have

$$
\|h(Y_k) - h(Y_k^{(m)})\|_p \leq K\|Y_k - Y_k^{(m)}\|^{\alpha} \leq K\|Y_k - Y_k^{(m)}\|_p^{\alpha},
$$

and thus $\{h(Y_k)\}$ is also weakly $\mathcal{M}$-dependent in $L^p$ with rate function $K\delta(\cdot)^\alpha$. 

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Note that (B) implies that for any fixed $m$ the sequence $\{Y_k^{(m)}, k \in \mathbb{Z}\}$ is an $m$-dependent process. Hence, sequences satisfying conditions (A) and (B) are approximable, in the $L^p$ sense, by $m$-dependent processes of any fixed order $m \geq 1$ with termwise approximation error $\delta(m)$. In other words, sequences in Definition 1 are close to $m$-dependent sequences, the value of $m$ depending on the required closeness, explaining the terminology. Since $\|Y_k\|_p \leq \|Y_k^{(m)}\|_p + \|Y_k - Y_k^{(m)}\|_p$, condition (A) implies that $E|Y_k|^p$ is finite. Using $L^p$-distance is convenient for our theorems, but, depending on the application, other distances can be used in part (A) of Definition 1. For example, defining (as usual) the $L^0$ norm of a random variable $X$ by

$$\|X\|_0 = \inf\{\varepsilon > 0 : P(|X| \geq \varepsilon) < \varepsilon\},$$

condition (A) could be replaced by

$$\|Y_k - Y_k^{(m)}\|_0 \leq \delta(m).$$

Such a definition requires no moment assumptions and turns out to provide a useful dependence measure for studying empirical processes (see [3]).

Trivially the previous definition covers $m$-dependent processes for any fixed $m$ (see also Section 3.1), but, in contrast to the very restrictive condition of $m$-dependence, weak $\mathcal{M}$-dependence holds for a huge class of stationary sequences, including those studied in Wu [49, 50] and Liu and Lin [28]. In the case when $\{Y_k, k \in \mathbb{Z}\}$ allows a Wiener–Rosenblatt representation

$$Y_k = f(\varepsilon_k, \varepsilon_{k-1}, \ldots), \quad k \in \mathbb{Z},$$

(6)

with an i.i.d. sequence $\{\varepsilon_k, k \in \mathbb{Z}\}$, weak $\mathcal{M}$-dependence is very close to Wu’s physical dependence condition in [49], except that we allow a larger freedom in choosing the approximating random variables $Y_k^{(m)}$, compared with the choice in [49, 50] via coupling. (For sufficient criteria for the representation (6), see Rosenblatt [40–42].) Note that instead of (6) we may also assume a two-sided representation

$$Y_k = f(\ldots, \varepsilon_{k-1}, \varepsilon_k, \varepsilon_{k+1}, \ldots), \quad k \in \mathbb{Z},$$

(7)

of $\{Y_k\}$. In case when $\{Y_k, k \in \mathbb{Z}\}$ allows the representation (7) with mixing $\{\varepsilon_k\}$, Definition 1 is a modified version of NED (see Section 3.2), a weak dependence condition which appeared already in Ibragimov [22] and has been brought forward in Billingsley [5] (see also [31, 32]). Later NED has been successfully used in the econometrics literature to establish weak dependence of dynamic time series models (see, e.g., [35]). In Section 3 we will discuss further the connection between weak $\mathcal{M}$-dependence with known weak dependence conditions. We stress that the definition of weak $\mathcal{M}$-dependence does not assume the representation (6) or (7), although it was
motivated by this case. The reason for using our more general definition is to illuminate the essential structural condition on \( \{Y_k\} \) required for our theorems. Extensions of our results for “classical” mixing conditions, like \( \alpha, \beta, \rho \) mixing and their variants will be given in a subsequent paper.

The main results of our paper are formulated in Section 2. In Section 3 we give several examples. Applications of the theorems can be found in Sections 4 and 5, while Section 6 contains the proofs of the main theorems.

2. Main theorems. We write \( a_n \ll b_n \) if \( \lim_{n \to \infty} |a_n/b_n| < \infty \).

**Theorem 1.** Let \( p > 2, \eta > 0 \) and let \( \{Y_k, k \in \mathbb{Z}\} \) be a centered stationary sequence, weakly \( \mathcal{M} \)-dependent in \( L^p \) with rate function
\[
\delta(m) \ll m^{-A},
\]
where
\[
A > \frac{p - 2}{2\eta} \left( 1 - \frac{1 + \eta}{p} \right) \vee 1, \quad (1 + \eta)/p < 1/2.
\]
Then the series
\[
\sigma^2 = \sum_{k \in \mathbb{Z}} EY_0Y_k
\]
is absolutely convergent, and \( \{Y_k, k \in \mathbb{Z}\} \) can be redefined on a new probability space together with two Wiener processes \( \{W_1(t), t \geq 0\} \) and \( \{W_2(t), t \geq 0\} \) such that
\[
\sum_{k=1}^{n} Y_k = W_1(s_n^2) + W_2(t_n^2) + O(n^{(1+\eta)/p}) \quad \text{a.s.,}
\]
where \( \{s_n\} \) and \( \{t_n\} \) are nondecreasing numerical sequences with
\[
s_n^2 \sim \sigma^2 n, \quad t_n^2 \sim cn^{\gamma}
\]
for some \( 0 < \gamma < 1, \ c > 0 \).

Note that for any fixed \( p > 2 \) and \( 0 < \eta < (p - 2)/2 \), condition (9) is satisfied if \( A \) is large enough, and thus Theorem 1 provides an a.s. invariance principle with remainder term close to the optimal remainder term \( o(n^{1/p}) \) in the Komlós–Major–Tusnády approximation.

It is natural to ask if \( W_1(s_n^2) \) in (11) can be replaced by \( W_1(\sigma^2 n) \), a fact that would simplify applications. The proof of the theorem yields an \( s_n \) with \( s_n^2 = \sigma^2 n + O(n^{1-\epsilon}) \) for some \( 0 < \epsilon < 1 \), but for \( A \) barely exceeding the lower bound in (9), the explicit value of \( \epsilon \) is very small. Thus replacing \( W_1(s_n^2) \) by \( W_1(\sigma^2 n) \) introduces an additional error term that ruins the error term \( O(n^{(1+\eta)/p}) \) in (11). The situation is similar to the Wiener approximation of partial sums of i.i.d. random variables with mean 0 and variance 1 when
we have (4) with a numerical sequence $\tau_n \sim n$, but in general (4) does not hold with $\tau_n = n$. (See Major [29, 30].) Note, however, that in our case the large difference between $s_n^2$ and $\sigma_n^2$ is a consequence of the method, and we do not claim that another construction cannot yield the approximation (11) with $s_n^2 = \sigma_n^2$. However, the presence of $s_n^2$ in (11) does not limit the applicability of our strong invariance principle: $s_n^2$ and $t_n^2$ are explicitly calculable nonrandom numbers and as we will see, applying limit theorems for $W_1(s_n^2)$ is as easy as for $W_1(\sigma_n^2)$.

As the proof of Theorem 1 will show, the sequences $\{s_n\}$ and $\{t_n\}$ in (11) have a complementary character. More precisely, there is a partition $N = G_1 \cup G_2$ (provided by the long and short blocks in a traditional blocking argument) and a representation

$$s_n^2 = \sum_{k=1}^{n} \sigma_k^2, \quad t_n^2 = \sum_{k=1}^{n} \tau_k^2 \quad (n = 1, 2, \ldots)$$

such that $\sigma_k^2$ converges to $\sigma^2$ on $G_1$ and equals 0 on $G_2$, and $\tau_k^2$ converges to $\sigma^2$ on $G_2$ and equals 0 on $G_1$. In particular,

$$\lim_{n \to \infty} (s_{n+1}^2 - s_n^2) = \lim_{n \to \infty} (t_{n+1}^2 - t_n^2) = \sigma^2,$$

and both limsup’s are equal to 0.

The numerical value of $\gamma$ in (12) plays no role in the applications in this paper, but for later applications we note that if

$$A > \frac{p - 2}{2\eta(1 - \varepsilon_0)^2} \left(1 - \frac{1 + \eta}{p}\right) \lor 1$$

for some $0 < \varepsilon_0 < 1$, then we can choose

$$\gamma = 1 - \varepsilon_0 \frac{2\eta(1 - \varepsilon_0)}{p - 2(1 + \eta\varepsilon_0)}.$$

As we already mentioned in the Introduction, the processes $W_1$ and $W_2$ are not independent. While for our applications this is not important, the following proposition might be useful for possible further applications.

**Proposition 1.** Under the assumptions of Theorem 1 we have

$$\text{Corr}(W_1(s_n), W_2(t_m)) \to 0 \quad \text{as } m, n \to \infty.$$

Our next theorem is the analogue of Theorem 1 in the case of an exponential decay in the dependence condition.

**Theorem 2.** Let $p > 2$ and let $\{Y_k, k \in \mathbb{Z}\}$ be a centered stationary sequence, weakly $\mathcal{M}$-dependent in $L^p$ with rate function

$$\delta(m) \ll \exp(-\varrho m), \quad \varrho > 0.$$
Then the series (10) is absolutely convergent, and the series \(\{Y_k, k \in \mathbb{Z}\}\) can be redefined on a new probability space together with two standard Wiener processes \(\{W_1(t), t \geq 0\}\) and \(\{W_2(t), t \geq 0\}\) such that
\[
\sum_{k=1}^{n} Y_k = W_1(s_n^2) + W_2(t_n^2) + O(n^{1/p} \log^2 n) \quad \text{a.s.,}
\]
where \(\{s_n\}\) and \(\{t_n\}\) are nondecreasing numerical sequences such that \(s_n^2 \sim \sigma^2 n\), \(t_n^2 \sim \sigma^2 n / \log n\) and (13) holds.

Like in Theorem 1, \(s_n^2 \sim \sigma^2 n\) can be sharpened to \(s_n^2 = \sigma^2 n + O(n / \log n)\); see the remarks after Theorem 1.

Using the law of the iterated logarithm for \(W_2\), relation (11) implies
\[
\sum_{k=1}^{n} Y_k = W_1(s_n^2) + O(n^{1/2 - \lambda}) \quad \text{a.s.}
\]
for some \(\lambda > 0\), which is the standard form of strong invariance principles. However, since \(\gamma\) in (12) is typically near to 1, the \(\lambda\) in (18) can be very small, and thus the effect of the very strong error term \(O(n^{(1+\eta)/p})\) in (11) is lost.

The proof of the strong approximation theorems in Wu [50] depends on martingale approximation, while Liu and Lin [28] use approximation of the partial sums of \(\{Y_k\}\) by partial sums of \(m\)-dependent r.v.’s. Our approach differs from both, using a direct approximation of separated block sums of \(\{Y_k\}\) by independent r.v.’s, an idea used earlier in [2–4, 21]. In this approach, the second Wiener process \(W_2\) is provided by the sum of short block sums. The question if one can get a remainder term near \(o(n^{1/p})\) in the simple (one-term) Wiener approximation for any \(p > 2\) remains open.

3. Examples of weakly \(\mathcal{M}\)-dependent processes. The classical approach to weak dependence, developed in the seminal papers of Rosenblatt [39] and Ibragimov [22], uses the strong mixing property and its variants like \(\beta\), \(\varrho\), \(\phi\) and \(\psi\) mixing, combined with a blocking technique to connect the partial sum behavior of \(\{Y_k\}\) with that of independent random variables. This method yields very sharp results (for a complete account of the classical theory see Bradley [7]), but verifying mixing conditions of the above type is not easy and even when they apply (e.g., for Markov processes), they typically require strong smoothness conditions on the process. For example, for the AR(1) process
\[
Y_k = \frac{1}{2} Y_{k-1} + \varepsilon_k
\]
with Bernoulli innovations, strong mixing fails to hold (cf. Andrews [1]). Recognizing this fact, an important line of research in probability theory in past years has been to find weak dependence conditions which are strong...
enough to imply satisfactory asymptotic results, but which are sufficiently
general to be satisfied in typical applications. Several conditions of this kind
have been found, in particular by the French school (see \[10–12, 16, 37, 38\]).
A different type of mixing conditions, the so-called physical and predictive
dependence measures, have been introduced by Wu \[49\] for stationary pro-
cesses \(\{Y_k\}\) admitting the representation (6) where \(\{\varepsilon_k, k \in \mathbb{Z}\}\) is an i.i.d.
sequence, and \(f : \mathbb{R}^N \to \mathbb{R}\) is a Borel-measurable function. These conditions
are particularly easy to handle, since they are defined in terms of the al-
gorithms which generate the process \(\{Y_k\}\). Weak \(\mathcal{M}\)-dependence, although
formally not requiring a representation of the form (6), is closely related to
Wu’s mixing conditions and works best for processes \(\{Y_k\}\) having a repre-
sentation (6) or its two-sided version (7). The examples below will clear up
the exact connection of our weak \(\mathcal{M}\)-dependence condition with the mixing
conditions in Wu \[49, 50\] and Liu and Lin \[28\].

3.1. \textit{m-dependent processes.} Definition 1 implies that \(\{Y_k, k \in \mathbb{Z}\}\) can be
approximated, for every \(m \geq 1\), by an \(m\)-dependent process with termwise \(L^p\)
error \(\delta(m)\). If \(\{Y_k, k \in \mathbb{Z}\}\) itself is \(m\)-dependent for some fixed \(m = m_0\) and
\(K := \sup_{k \in \mathbb{Z}} \|Y_k\|_p < \infty\), then Definition 1 is satisfied with
\[
\delta(j) = \begin{cases} K, & \text{if } j < m_0, \\ 0, & \text{if } j \geq m_0, \end{cases}
\]
and \(Y_k^{(n)} = 0\) if \(n < m_0\) and \(Y_k^{(n)} = Y_k\) if \(n \geq m_0\). In other words, \(m\)-dependent
sequences with uniformly bounded \(L^p\) norms are weakly \(\mathcal{M}\)-dependent with
the above parameters. It is worth mentioning that \(m\)-dependent processes
in general do not have the representation (7) (see, e.g., \[8, 13\]).

3.2. \textit{NED processes.} Under (7) our condition can be directly compared
to NED. We recall:

**Definition 2 (NED).** A sequence \(\{Y_k, k \in \mathbb{Z}\}\) having representation (7)
is called \textit{NED} over \(\{\varepsilon_k\}\) under \(L^p\)-norm with rate function \(\delta(\cdot)\) if for any
\(k \in \mathbb{Z}, m \geq 1\),
\[
\|Y_k - E[Y_k|\mathcal{F}_{k-m}^{k+m}]\|_p \leq \delta(m),
\]
where \(\mathcal{F}_{k-m}^{k+m}\) is the \(\sigma\)-algebra generated by \(\varepsilon_{k-m}, \ldots, \varepsilon_{k+m}\).

Clearly, if \(\{\varepsilon_k\}\) is an independent sequence, then \(Y_k^{(m)} = E[Y_k|\mathcal{F}_{k-m}^{k+m}]\) satisfies
(B) of Definition 1. Hence if \(\{Y_k\}\) is NED over \(\{\varepsilon_k\}\) in \(L^p\)-norm with
rate function \(\delta(\cdot)\) where \(\{\varepsilon_k\}\) is an independent sequence, then \(\{Y_k\}\) is weakly
\(\mathcal{M}\)-dependent with the same \(p, \delta(\cdot)\).

As our examples below will show, for weakly \(\mathcal{M}\)-dependent sequences the
construction for \(Y_k^{(m)}\) is not restricted to \(E[Y_k|\mathcal{F}_{k-m}^{k+m}]\), but is often more
conveniently established by truncation or coupling methods.
3.3. Linear processes. Let \( Y_k = \sum_{j=-\infty}^{\infty} a_j \varepsilon_{k-j} \) with the i.i.d. innovations \( \{\varepsilon_j, j \in \mathbb{Z}\} \). If \( a_j = 0 \) for \( j < 0 \), then the sequence \( \{Y_k, k \in \mathbb{Z}\} \) is causal. Liu and Lin \[28\] and Wang, Lin and Gulati \[48\] studied strong approximations of the partial sums with Gaussian processes (in the short- and long-memory cases).

We define \( Y_k^{(m)} \) as

\[
Y_k^{(m)} = \sum_{j=-[m/2]}^{[m/2]} a_j \varepsilon_{k-j}.
\]

This directly ensures that condition (B) holds. To verify condition (8) we will assume that \( E|\varepsilon_0|^p < \infty \) for some \( p > 2 \) as well as \( |a_j| \ll |j|^{-(A+1)} \) \((j \to \infty)\). Then we get, using the Minkowski inequality,

\[
\|Y_k - Y_k^{(m)}\|_p = \left\| \sum_{|j| > m/2} a_j \varepsilon_{k-j} \right\|_p \\
\leq \sum_{|j| > m/2} \|a_j \varepsilon_{k-j}\|_p \\
= (E|\varepsilon_0|^p)^{1/p} \sum_{|j| > m/2} |a_j| \ll m^{-A}.
\]

Thus if \( A \) is large enough, Theorem 1 applies. Obviously if \( |a_j| \ll \rho|j| \) with some \( 0 < \rho < 1 \), then (16) holds, and Theorem 2 applies.

3.4. Nonlinear time series. Let the time series \( \{Y_k, k \in \mathbb{Z}\} \) be defined by the stochastic recurrence equation

\[
Y_k = G(Y_{k-1}, \varepsilon_k),
\]

where \( G \) is a measurable function, and \( \{\varepsilon_k, k \in \mathbb{Z}\} \) is an i.i.d. sequence. For example, ARCH(1) processes (see, e.g., Engle \[19\]) which play an important role in the econometrics literature, are included in this setting. Sufficient conditions for the existence of a stationary solution of (19) can be found in Diaconis and Freedman \[14\]. Note that iterating (19) yields \( Y_k = f(\ldots, \varepsilon_{k-1}, \varepsilon_k) \) for some measurable function \( f \). This suggests defining the approximating random variables \( Y_k^{(m)} \) as \( Y_k^{(m)} = f(\ldots, 0, 0, \varepsilon_{k-m}, \ldots, \varepsilon_k) \). Note, however, that this definition does not guarantee the convergence and thus the existence of \( Y_k^{(m)} \). The coupling used by Wu \[49\], avoids this problem by defining

\[
Y_k^{(m)} = f(\ldots, \varepsilon_{k-m-2}^{(k)}, \varepsilon_{k-m-1}^{(k)}, \varepsilon_{k-m}, \ldots, \varepsilon_k),
\]

where \( \{\varepsilon_l^{(l)}, k \in \mathbb{Z}\}, l = 1, 2, \ldots, \) are i.i.d. sequences with the same distribution as \( \{\varepsilon_k, k \in \mathbb{Z}\} \) which are independent of each other and of the \( \{\varepsilon_k, k \in \mathbb{Z}\} \). These random variables satisfy condition (B). Results from Wu
Shao [52] show that under some simple technical assumption on $G$,
\[ \|Y_k - Y_k^{(m)}\|_p \ll \exp(-\rho m) \]
holds with some $p > 0$ and $\rho > 0$. Thus for $p > 2$, Theorem 2 applies.

3.5. Augmented GARCH sequences. Augmented GARCH sequences were introduced by Duan [17] and turned out to be very useful in applications in macroeconomics and finance. The model is quite general and many popular processes are included in its framework. Among others the well-known GARCH [6], AGARCH [15] and EGARCH model [34] are covered. We consider the special case of augmented GARCH(1,1) sequences, that is, sequences $\{Y_k, k \in \mathbb{Z}\}$ defined by
\[ Y_k = \sigma_k \varepsilon_k, \]
(20)
where the conditional variance $\sigma_k^2$ is given by
\[ \Lambda(\sigma_k^2) = c(\varepsilon_{k-1}) \Lambda(\sigma_{k-1}^2) + g(\varepsilon_{k-1}). \]
(21)
Here $\{\varepsilon_k, k \in \mathbb{Z}\}$ is a sequence of i.i.d. errors, and $\Lambda(x)$, $c(x)$ and $g(x)$ are real-valued measurable functions. To solve (21) for $\sigma_k^2$ one usually assumes that $\Lambda^{-1}(x)$ exists. Necessary and sufficient conditions for the existence of a strictly stationary solution of (20) and (21) were given by Duan [17] and Aue, Berkes and Horváth [2]. Under some technical conditions stated in Hörmann [21] (Lemmas 1, 2 and Remark 2) one can show that augmented GARCH sequences are weakly $M$-dependent in $L^p$-norm with exponential rate.

Note that the above models have short memory; long memory models (see, e.g., [20]) have completely different properties.

3.6. Linear processes with dependent innovations. Linear processes $Z_k = \sum_{j=-\infty}^{\infty} a_j Y_{k-j}$ with dependent innovations $\{Y_k\}$ have obtained considerable interest in the financial literature. A common example are autoregressive (AR) processes with augmented GARCH innovations (see, e.g., [27]).

Assume that $\{Y_k, k \in \mathbb{Z}\}$ is weakly $M$-dependent in $L^p$ with rate function $\delta(\cdot)$. In combination with the results of Section 3.3 one can easily obtain conditions on $\delta$ assuring that the linear process $\{Z_k, k \in \mathbb{Z}\}$ defined above is also weakly $M$-dependent in $L^p$-norm with a rate function $\delta^*$ depending on $(a_j)$ and $\delta$.

Strong approximation results for linear processes with dependent errors were also obtained by Wu and Min [51].

3.7. Ergodic sums. Let $f$ be a real measurable function with period 1 such that $\int_{0}^{1} f(\omega) \, d\omega = 0$ and $\int_{0}^{1} |f(\omega)|^p \, d\omega < \infty$ for some $p > 2$. Set
\[ S_n(\omega) = \sum_{k=1}^{n} f(2^k \omega), \quad \omega \in [0, 1), \]
SPLIT INVARIANCE PRINCIPLES

and \( B_n^2 = \int_0^1 S_n^2(\omega) \, d\omega \). Then \( S_n \) defines a partial sum process on the probability space \(((0,1), \mathcal{B}_{[0,1]}, \lambda_{[0,1]})\), where \( \mathcal{B}_{[0,1]} \) and \( \lambda_{[0,1]} \) are the Borel \( \sigma \)-algebra and Lebesgue measure on \([0,1)\). The strong law of large numbers for \( f(2^k \omega) \) is a consequence of the ergodic theorem, for central and functional central limit theorems see Kac [24], Ibragimov [23] and Billingsley [5].

Let \( Y_k(\omega) = f(2^k \omega) \), and define the random variable \( \varepsilon_k(\omega) \) to be equal to the \( k \)th digit in the binary expansion of \( \omega \). Ambiguity can be avoided by the convention to take terminating expansions whenever possible. Then \( \{\varepsilon_k\} \) is an i.i.d. sequence, and we have \( \varepsilon_k = \pm 1 \), each with probability \( 1/2 \). This gives the representation

\[
Y_k = f \left( \sum_{j=1}^{\infty} \varepsilon_{k+j} 2^{-j} \right) = g(\varepsilon_{k+1}, \varepsilon_{k+2}, \ldots).
\]

We can now make use of the coupling method described in Section 3.4 and the approximations

\[
Y_k^{(m)} = g(\varepsilon_{k+1}, \varepsilon_{k+2}, \ldots, \varepsilon_{k+m}, \varepsilon_{k+m+1}, \varepsilon_{k+m+2}, \ldots).
\]

Changing for some \( \omega \in [0,1) \) the digits \( \varepsilon_k(\omega) \) for \( k > m \) will give an \( \omega' \) with \( |\omega - \omega'| \leq 2^{-m} \). If \( f \) is Lipschitz continuous of some order \( \gamma \), then we have

\[
|Y_k - Y_k^{(m)}| = O(2^{-\gamma m}),
\]

and thus for any \( p \geq 1 \) \( \{Y_k\} \) is weakly \( \mathcal{M} \)-dependent in \( L^p \)-norm with an exponentially decaying rate function.

4. Increments of stochastic processes. For arbitrary \( \lambda > 0 \), relation (18) has many useful applications in probability and statistics. For example, it implies a large class of limit theorems on CLT and LIL behavior and for various other functionals of weakly dependent sequences. However, many refined limit theorems for partial sums require a remainder term better than \( O(n^{1/4}) \), and no existing method for dependent sequences provides such a remainder term. The purpose of the next two sections is to show how to deal with such limit theorems via our approximation results in Section 2.

Let \( \{Y_k, k \in \mathbb{Z}\} \) be a stationary random sequence, and let \( 0 < a_n \leq n \) be a nondecreasing sequence of real numbers. In this section, we investigate the order of magnitude of

\[
\max_{1 \leq k \leq n-a_n} \max_{1 \leq \ell \leq a_n} \left| \sum_{j=k+1}^{k+\ell} Y_j \right|.
\]

Such results have been obtained by Csörgő and Révész [9] for i.i.d. sequences and the Wiener process. In particular, they obtained the following result ([9], Theorem 1.2.1).
Theorem 3. Let \( \{a_T, T \geq 0\} \) be a positive nondecreasing function satisfying:

(a) \( 0 < a_T \leq T \);
(b) \( T/a_T \) is nondecreasing.

Set

\[
\beta_T = \left( 2a_T \left[ \log \frac{T}{a_T} + \log \log T \right] \right)^{-1/2}.
\]

Then

\[
\varlimsup_{T \to \infty} \max_{0 \leq t \leq T-a_T} \max_{0 \leq s \leq a_T} \beta_T |W(t + s) - W(t)| = 1.
\]

Using strong invariance, a similar result can be obtained for partial sums of i.i.d. random variables under suitable moment conditions (see [9], pages 115–118). For slowly growing \( a_T \), this requires a very good remainder term in the Wiener approximation of partial sums, using the full power of the Komlós–Major–Tusnády theorems. As an application of our main theorems in Section 2, we now extend Theorem 3 for dependent stationary processes. To simplify the formulation and to clarify the connection between the remainder term in our approximation theorems in Section 2 and the increment problem, we introduce the following assumption.

Assumption 1. Let \( \{Y_k\} \) be a random sequence which can be redefined on a new probability space together with two standard Wiener processes \( \{W_1(t), t \geq 0\} \) and \( \{W_2(t), t \geq 0\} \) such that

\[
\sum_{k=1}^{n} Y_k = W_1(s_n^2) + W_2(t_n^2) + O(E_n) \quad \text{a.s.,}
\]

where \( \{E_n\} \) is some given sequence and \( \{s_n^2\} \) and \( \{t_n^2\} \) are nondecreasing sequences satisfying

\[
s_n^2 \sim \sigma^2 n, \quad t_n^2 = o(n), \quad \lim_{k \to \infty} (s_{k+1}^2 - s_k^2) = \lim_{k \to \infty} (t_{k+1}^2 - t_k^2) = \sigma^2.
\]

We will prove the following result.

Theorem 4. Let \( \{Y_k\} \) be a sequence of random variables satisfying Assumption 1 and put \( S_n = \sum_{k=1}^{n} Y_k \). Let \( a_T \) be a positive nondecreasing function such that:

(a) \( 0 < a_T \leq T \);
(b) \( T/a_T \) is nondecreasing;
(c) \( a_T \) is regularly varying at \( \infty \) with index \( q \in (0, 1] \).
Let $\beta_T$ be defined by (22). Then under the condition

$$\beta_T E_T = o(1)$$

we have

$$\lim_{n \to \infty} \max_{1 \leq k \leq n-a_n} \max_{1 \leq \ell \leq a_n} \beta_n |S_{k+\ell} - S_k| = \sigma^2.$$  \hfill (26)

Given a function $a_T$ and a weakly $\mathcal{M}$-dependent sequence $\{Y_k\}$ with parameters $p, \delta(\cdot)$, we can compute, using Theorem 4, a rate of decrease for $\delta(\cdot)$ and a value for $p > 2$ such that the fluctuation result (26) holds. For example, if $a_T = T^\alpha$, $0 < \alpha < 1$, then (26) holds if $p > 4/\alpha$ and $\delta(m) \ll m^{-p/2}$.

We note that for i.i.d. observations only assumptions (a) and (b) are required. It remains open whether a more general version of our Theorem 4 which does not require assumption (c) can be proved.

Recently Zholud [53] obtained a distributional version of Theorem 3 by showing that the functional

$$\max_{0 \leq t \leq T-a_T} \max_{0 \leq s \leq a_T} (W(t+s) - W(t))$$

converges weakly, suitably centered and normalized, to the extremal distribution with distribution function $e^{-e^{-x}}$. Using this fact and our a.s. invariance principles, a distributional version of Theorem 4 can be obtained easily. Since the argument is similar to that for (26), we omit the details.

Let $|A|$ be the cardinality of a set $A$. For the proof of Theorem 4 we need the following simple lemma.

**Lemma 1.** Assume that $\{d_k, k \geq 1\}$ is a nonincreasing sequence of positive numbers such that $\sum_{k=1}^{\infty} d_k = \infty$. Let $A \subset \mathbb{N}$ have positive density, that is,

$$\liminf_{n \to \infty} |A \cap \{1, \ldots, n\}|/n > 0.$$  

Then $\sum_{k=1}^{\infty} d_k I\{k \in A\} = \infty$.

**Proof.** First note that by our assumption we have $\sum_{k=1}^{n} I\{k \in A\} \geq \mu n$ for some $\mu > 0$ as long as $n \geq n_0$. Using Abel summation we can write

$$\sum_{k=1}^{n} d_k = nd_n + \sum_{k=1}^{n-1} k(d_k - d_{k+1}).$$

Hence, by our assumptions

$$nd_n + \sum_{k=n_0}^{n-1} k(d_k - d_{k+1}) \to \infty \quad (n \to \infty).$$
From \(d_k - d_{k+1} \geq 0\) it follows (again using the Abel summation) that for \(n \geq n_0\)
\[
\sum_{k=1}^{n} d_k I\{k \in A\} = d_n \sum_{k=1}^{n} I\{k \in A\} + \sum_{k=1}^{n-1} (d_k - d_{k+1}) \sum_{j=1}^{k} I\{j \in A\} \\
\geq \mu \left( n d_n + \sum_{k=n_0}^{n-1} k(d_k - d_{k+1}) \right) \to \infty \quad (n \to \infty).
\]

**Proof of Theorem 4.** For the sake of simplicity we carry out the proof for \(\sigma = 1\). From (23) and the triangular inequality we infer that
\[
\lim_{n \to \infty} \max_{1 \leq k \leq n-a_n} \max_{1 \leq \ell \leq a_n} \beta_n |S_{k+\ell} - S_k| \\
\leq \lim_{n \to \infty} \max_{1 \leq k \leq n-a_n} \max_{1 \leq \ell \leq a_n} \beta_n |W_1(s_{k+\ell}^2) - W_1(s_k^2)| \\
+ \lim_{n \to \infty} \max_{1 \leq k \leq n-a_n} \max_{1 \leq \ell \leq a_n} \beta_n |W_2(t_{k+\ell}^2) - W_2(t_k^2)| \\
+ \lim_{n \to \infty} \beta_n O(E_n) \\
= A_1 + A_2 + A_3.
\]

By (25) \(A_3 = 0\). Since \(a_n \to \infty\) [this is implicit in (c)], we conclude from (24) that for any \(\varepsilon > 0\) some \(n_0\) exists, such that for all \(n \geq n_0\)
\[
\sup_{k \geq 1} \{s_{k+a_n}^2 - s_k^2\} \leq (1 + \varepsilon)a_n \quad \text{and} \quad s_n^2 \leq (1 + \varepsilon)n.
\]
Set \(T = (1 + \varepsilon)n\), and define \(a_{T,\varepsilon} = (1 + \varepsilon)a_{T/(1+\varepsilon)}\). Then \(a_{T,\varepsilon}\) satisfies (a) and (b) and for \(n \geq n_0\) we have
\[
\max_{1 \leq k \leq n-a_n} \max_{1 \leq \ell \leq a_n} \beta_n |W_1(s_{k+\ell}^2) - W_1(s_k^2)| \\
\leq \sup_{0 \leq t \leq s_{n-a_n}^2} \sup_{0 \leq s \leq (1+\varepsilon)a_n} \beta_n |W_1(t + s) - W_1(t)| \\
\leq \sup_{0 \leq t \leq T - a_{T,\varepsilon}} \sup_{0 \leq s \leq a_{T,\varepsilon}} \beta_{T/(1+\varepsilon)} |W_1(t + s) - W_1(t)|.
\]

Let
\[
\beta_{T,\varepsilon} = \left(2a_{T,\varepsilon} \left[\log \frac{T}{a_{T,\varepsilon}} + \log \log T\right]\right)^{-1/2}.
\]

By application of Theorem 1.2.1 in Csörgö and Révész [9] [which requires (a) and (b)] we get
\[
\lim_{T \to \infty} \sup_{0 \leq t \leq T - a_{T,\varepsilon}} \sup_{0 \leq s \leq a_{T,\varepsilon}} \beta_{T,\varepsilon} |W_1(t + s) - W_1(t)| = 1 \quad \text{a.s.}
\]
Since $\lim_{T \to \infty} \beta_T/\beta_T = (1 + \varepsilon)^{1/2}$, and $\varepsilon$ can be chosen arbitrarily small, we have shown that $A_1 \leq 1$ a.s.

It is not surprising that due to (24) similar arguments will lead to $A_2 = 0$ a.s.

The proof will be completed if we show that $A_1 \geq 1$. Let $\{n_k\}$ be a nondecreasing sequence of integers with $n_k \to \infty$. By (23), the triangular inequality and $A_2 = A_3 = 0$ we obtain

$$
\lim_{n \to \infty} \max_{1 \leq k \leq n-a_n} \max_{1 \leq \ell \leq a_n} \beta_n |S_{k+\ell} - S_k|
\geq \lim_{n \to \infty} \max_{1 \leq k \leq n-a_n} \max_{1 \leq \ell \leq a_n} \beta_n |W_1(s^2_{k+\ell}) - W_1(s^2_k)|
\geq \lim_{k \to \infty} \beta_{n_k} |W_1(s^2_{n_k}) - W_1(s^2_{n_k-a_{n_k}})|.
$$

We now proceed similarly as in Csörgő and Révesz [9] for the proof of Step 2 of their Theorem 1.2.1. We will distinguish between the cases $\lim a_T/T = \rho$ with $\rho < 1$ and $\rho = 1$. Since both times we can use the same conceptual idea, we shall treat here only $\rho < 1$.

Set $n_1 = 1$. Given $n_k$, define $n_{k+1}$ such that $n_{k+1} - a_{n_{k+1}} = n_k$. This equation will, in general, have no integer solutions, but for the sake of simplicity we assume that $(n_k)$ and $(a_{n_k})$ are $\mathbb{Z}$-valued. Since $(s^2_{n_k})$ is nondecreasing, we conclude that the increments $\Delta(k) = W_1(s^2_{n_k}) - W_1(s^2_{n_k-a_{n_k}})$ are independent.

By the second Borel–Cantelli lemma it suffices to show now that

$$
(27) \quad \sum_{k=1}^{\infty} P(\beta_{n_k} |\Delta(k)| \geq 1 - \varepsilon) = \infty \quad \text{for all } \varepsilon > 0.
$$

For all large enough $k \in \mathbb{N}$ for which $s^2_{n_k} - s^2_{n_{k-1}} \geq (1 - \varepsilon/2)a_{n_k}$, the estimates in [9] give

$$
P(\beta_{n_k} |\Delta(k)| \geq 1 - \varepsilon) \geq \left(\frac{a_{n_k}}{n_k \log n_k}\right)^{1-\varepsilon}.
$$

It is also shown in [9] that $\sum_{k=1}^{\infty} (\frac{a_{n_k}}{n_k \log n_k})^{1-\varepsilon} = \infty$. Thus, in view of Lemma 1 it remains to show that $A = \{k \geq 1 | s^2_{n_k} - s^2_{n_{k-1}} \geq (1 - \varepsilon/2)a_{n_k}\}$ has a positive density. By (24) we have

$$
(s^2_{n_k} - s^2_{n_{k-1}})/n_k = \sum_{j=2}^{k} (s^2_{n_j} - s^2_{n_{j-1}})/n_k
\leq C_0 \sum_{2 \leq j \leq k} (n_j - n_{j-1})/n_k + \sum_{2 \leq j \leq k} (1 - \varepsilon/2)(n_j - n_{j-1})/n_k
$$

Since $\lim_{T \to \infty} \beta_T/\beta_T = (1 + \varepsilon)^{1/2}$, and $\varepsilon$ can be chosen arbitrarily small, we have shown that $A_1 \leq 1$ a.s.
\[ \leq C_0 \sum_{2 \leq j \leq k, j \in A} (n_j - n_{j-1})/n_k + (1 - n_1/n_k)(1 - \epsilon/2) \]

for some \( C_0 > 0 \) which is independent of \( k \). Now if \( A \) had density zero, the limsup of the right-hand side of the last relation would be \( 1 - \epsilon/2 \). This can be easily proved, using that \( (n_j - n_{j-1}) \) is regularly varying by assumption (c). The liminf of the left-hand side above is 1. Thus \( A \) must have positive density and the proof is complete. \( \square \)

5. Change-point tests with an epidemic alternative. In this section we apply our invariance principles to a change-point problem. Let \( \{Y_k, k \in \mathbb{Z}\} \) be a zero mean process. Further let \( X_k = Y_k + \mu_k \), where \( \mu_k, k \in \mathbb{Z} \), are unknown constants. We want to test the hypothesis

\[(H_0) \quad \mu_1 = \mu_2 = \cdots = \mu_n = \mu\]

against the “epidemic alternative”

\[(H_A) \quad \text{There exist } 1 \leq m_1 < m_2 \leq n \text{ such that } \mu_k = \mu \text{ for } k \in \{1, \ldots, n\} \setminus \{m_1 + 1, \ldots, m_2\} \text{ and } \mu_k = \mu + \Delta \text{ if } k \in \{m_1 + 1, \ldots, m_2\}.\]

It should be noted that the variables \( m_1, m_2 \) and \( \Delta \) may depend on the sample size \( n \). As it is common in the change-point literature, this dependence is suppressed in the notation.

Without loss of generality we assume that \( \sigma = 1 \). To detect a possible epidemic change it is natural to compare the increments of the process to a proportion of the total sum. More specifically, assume for the moment that \( X_k \) are independent and that we know when the epidemic starts and ends. Set \( S_k = X_1 + \cdots + X_k \). Then by the law of large numbers \( I(m_1, m_2) = |S_{m_2} - S_{m_1} - (m_2 - m_1)S_n/n| > m_2 - m_1 \). If no change occurs, however, by the central limit theorem \( I(m_1, m_2) = O_P(\sqrt{m_2 - m_1}) \). In general we do not know \( m_1 \) and \( m_2 \). Thus, a natural test statistic is

\[ \max_{1 \leq i < j \leq n} |S_j - S_i - (j - i)S_n/n|. \]

Clearly we are required to normalize the above test statistic appropriately. Following Račkauskas and Suquet [36] we define

\[ UT(n, \alpha) = n^{-1/2} \max_{1 \leq i < j \leq n} \frac{|S_j - S_i - (j - i)S_n/n|}{\left[ ((j - i)/n)(1 - (j - i)/n) \right]^\alpha} \]

with \( 0 < \alpha < 1/2 \). As we will see below, the parameter \( \alpha \) plays an important role. The closer \( \alpha \) is to \( 1/2 \), the “shorter” epidemics can be detected with this test. The price, however, is that in order to obtain the limiting law under \( (H_0) \) with “large” \( \alpha \) (close to \( 1/2 \)) requires a.s. invariance principles with error \( n^{\epsilon} \), \( \epsilon \) close to zero. Choosing \( \alpha \geq 1/2 \) would result in a degenerate limiting distribution under \( (H_0) \).
PROPOSITION 2 [Asymptotics under \((H_0)\)]. If the stationary sequence \(\{Y_k, k \in \mathbb{Z}\}\) satisfies Assumption 1 with \(E_n = o(n^{1/2-\alpha})\) and \((H_0)\) holds, then
\[
\sigma^{-1}UI(n, \alpha) \xrightarrow{D} \sup_{0 < s < t < 1} \frac{|B(t) - B(s)|}{(t-s)(1-(t-s))^{\alpha}},
\]
where \(\{B(t), t \in [0,1]\}\) is a Brownian bridge.

PROOF. Using (23) and assuming for simplicity that \(\sigma = 1\), we obtain
\[
UI(n, \alpha) \leq n^{-1/2} \max_{1 \leq i < j \leq n} \frac{|W_1(s_j^2) - W_1(s_i^2) - (j-i)W_1(s_n^2)/n|}{((j-i)/n)(1-(j-i)/n)^\alpha} + n^{-1/2} \max_{1 \leq i < j \leq n} \frac{|W_2(t_j^2) - W_2(t_i^2) - (j-i)W_2(t_n^2)/n|}{((j-i)/n)(1-(j-i)/n)^\alpha}
\]
\[+ O(n^{-1/2+\alpha}E_n) = n^{-1/2} \max_{1 \leq i < j \leq n} T_{i,j}^{(1)} + n^{-1/2} \max_{1 \leq i < j \leq n} T_{i,j}^{(2)} + o(1). \]

It is easy to see that \(n^{-1/2} \max_{1 \leq i < j \leq n} T_{i,j}^{(2)}\) tends to zero. Since we can get a similar lower bound for \(UI(n, \alpha)\), we have
\[
UI(n, \alpha) = n^{-1/2} \max_{(i,j) \in M_n} T_{i,j}^{(1)} + o_P(1),
\]
where \(M_n = \{(i,j)|1 \leq i < j \leq n\}\). Let us partition \(M_n\) into
\[
M_{1,n} = \{(i,j)|1 \leq i < j \leq n;n \gamma_n < j - i < n(1-\gamma_n)\},
\]
\[
M_{2,n} = \{(i,j)|1 \leq i < j \leq n;n \gamma_n \geq j - i\}
\]
and
\[
M_{3,n} = \{(i,j)|1 \leq i < j \leq n;j - i \geq n(1-\gamma_n)\},
\]
where \(\gamma_n \to 0\) will be defined later. By our assumptions on the sequence \(\{s_j^2\}\) there exists a \(\tau > 0\) such that \(s_j^2 - s_i^2 \leq \tau(j-i)\) for all \(1 \leq i \leq j\) and that \(s_n^2 \leq (2 - \tau \gamma_n)n\) if \(n \geq n_0\). We have for large enough \(n\)
\[
n^{-1/2} \max_{(i,j) \in M_{2,n}} T_{i,j}^{(1)}
\]
\[\leq 2n^{\alpha - 1/2} \max_{(i,j) \in M_{2,n}} \left\{ \frac{|W_1(s_j^2) - W_1(s_i^2)|}{(j-i)^\alpha} \right\} + 2n^{-1/2} \tau n^{-1} |W_1(s_n^2)|
\]
\[\leq 2n^{\alpha - 1/2} \max_{1 \leq h \leq n \gamma_n} \sup_{0 \leq t \leq (2 - \tau \gamma_n)n} \sup_{0 \leq s \leq \tau h} \left\{ \frac{|W_1(t+s) - W_1(t)|}{h^\alpha} \right\} + o_P(1). \]
For arbitrary $\epsilon > 0$ we get by Lemma 1.2.1 in Csörgő and Révész [9] that there is a constant $C$ which is independent of $n$ and $\epsilon$ such that

$$P\left(\max_{1 \leq h \leq n^{\gamma_0}} \sup_{0 \leq t \leq (2-\epsilon)n} \sup_{0 \leq s \leq \tau h} \left\{ \frac{|W_1(t + s) - W_1(t)|}{h^\alpha} \right\} > \epsilon n^{1/2 - \alpha} \right)$$

$$\leq \sum_{h=1}^{n^{\delta_0}} P\left( \sup_{0 \leq t \leq 2n - \tau h} \sup_{0 \leq s \leq \tau h} |W_1(t + s) - W_1(t)| > \epsilon h^{1/2} (n/h)^{1/2 - \alpha} \right)$$

$$\leq \sum_{h=1}^{n^{\delta_0}} C n \epsilon e^{-\left(\epsilon^2/3\right)(n/h)^{1-2\alpha}} \to 0 \quad (n \to \infty).$$

Hence $n^{-1/2} \max_{(i,j) \in \mathcal{M}_{2,n}} T_{i,j}^{(1)} = o_P(1)$. In the same fashion one can show that $n^{-1/2} \max_{(i,j) \in \mathcal{M}_{3,n}} T_{i,j}^{(1)} = o_P(1)$. Therefore

$$UI(n, \alpha) = n^{-1/2} \max_{(i,j) \in \mathcal{M}_{1,n}} T_{i,j}^{(1)} + o_P(1).$$

Some further basic estimates give

$$n^{-1/2} \max_{(i,j) \in \mathcal{M}_{1,n}} T_{i,j}^{(1)} = n^{-1/2} \max_{(i,j) \in \mathcal{M}_{1,n}} \frac{|W_1(j) - W_1(i) - (j-i)W_1(n)/n|}{[(j-i)/n(1-(j-i)/n)]^\alpha}$$

$$+ O\left( n^{-1/2} \max_{1 \leq t \leq n} |W_1(i) - W_1(s_i^2)| \right).$$

Since $s_i^2 \sim n$ there is a null sequence $\{\epsilon_n\}$ such that $\max_{1 \leq i \leq n} |i - s_i^2| \leq \epsilon_n n$. Hence

$$\max_{1 \leq i \leq n} |W_1(i) - W_1(s_i^2)| \leq \sup_{0 \leq t \leq n} \sup_{0 \leq s \leq 2\epsilon_n n} |W_1(t + s) - W_1(t)|,$$

Setting $\gamma_n = \epsilon_n$ and applying again Lemma 1.2.1 in [9] it can be seen that

$$\sup_{0 \leq t \leq n} \sup_{0 \leq s \leq 2\epsilon_n n} |W_1(t + s) - W_1(t)| = o_P(n^{1/2 \gamma_n^\alpha}).$$

Consequently

$$UI(n, \alpha) = n^{-1/2} \max_{(i,j) \in \mathcal{M}_{1,n}} \frac{|W_1(j) - W_1(i) - (j-i)W_1(n)/n|}{[(j-i)/n(1-(j-i)/n)]^\alpha}$$

$$+ o_P(1).$$

Since the line of argumentation is very similar to what we have shown before, we note now without proof that $\mathcal{M}_{1,n}$ in the right-hand side of (28) can be replaced by $\mathcal{M}_n$. The rest of the proof of Proposition 2 is standard. □
The next proposition shows that this test is consistent. Let \( \ell = m_2 - m_1 \) denote the length of the epidemic.

**Proposition 3 [Asymptotics under \((H_A)\)].** Let \( \{Y_k, k \in \mathbb{Z}\} \) be a mean zero process, weakly \( \mathcal{M} \)-dependent in \( L^p \) with \( p \geq 2 \) and \( \delta(\cdot) \) satisfying

\[
\sum_{m \geq 1} \delta(m) < \infty.
\]

Let \( X_k = Y_k + \mu_k, k \in \mathbb{Z} \). Assume that \((H_A)\) holds and that

\[
\lim_{n \to \infty} \frac{(\ell(n - \ell))^{1-\alpha}}{n^{3/2-2\alpha}} |\Delta| = \infty.
\]

Then \( \text{UI}(n, \alpha) \xrightarrow{p} \infty \).

**Proof.** Under the alternative hypothesis \((H_A)\) we have \( X_k = Y_k + \mu \) for \( k \in \{1, \ldots, n\} \setminus \{m_1 + 1, \ldots, m_2\} \) and \( X_k = Y_k + \mu + \Delta \) for \( k \in \{m_1 + 1, \ldots, m_2\} \). To find a lower bound for \( \text{UI}(n, \alpha) \) we study the numerator of the test statistic corresponding to the true epidemic. Thus we look at

\[
S_{m_2} - S_{m_1} - S_n(m_2/n - m_1/n)
\]

\[
= (1 - \ell/n)(S_{m_2} - S_{m_1}) - (\ell/n)(S_n - (S_{m_2} - S_{m_1}))
\]

\[
= \frac{\ell(n - \ell)}{n} \Delta + (1 - \ell/n) \sum_{j=m_1+1}^{m_2} Y_j - (\ell/n) \left( \sum_{j=1}^{m_1} Y_j + \sum_{j=m_2+1}^{n} Y_j \right)
\]

\[
= \frac{\ell(n - \ell)}{n} \Delta + R_n.
\]

With the help of the moment inequality stated in Proposition 4 below we get

\[
\text{Var}(n^{-1/2}R_n) = O((1 - \ell/n)^2(\ell/n) + (\ell/n)^2(1 - \ell/n))
\]

\[
+ 2(\ell/n)^{3/2}(1 - \ell/n)^{3/2}
\]

\[
= O((\ell/n)(1 - \ell/n)),
\]

and thus \( n^{-1/2}R_n = O_P((\ell/n)^{1/2}(1 - \ell/n)^{1/2}) \). Thus we have shown that

\[
\text{UI}(n, \alpha) \geq n^{1/2}((\ell/n)(1 - \ell/n))^{1-\alpha} |\Delta|
\]

\[
- O_P(((\ell/n)(1 - \ell/n))^{1/2-\alpha})
\]

\[
= \frac{(\ell(n - \ell))^{1-\alpha}}{n^{3/2-2\alpha}} |\Delta|
\]

\[
- O_P(((\ell/n)(1 - \ell/n))^{1/2-\alpha}).
\]
To conclude the proof we note that \( \lim_{n \to \infty} (\ell/n)(1 - \ell/n)^{1/2 - \alpha} = 0 \) if \( \ell = o(n) \) [or \( n - \ell = o(n) \), resp.] and \((\ell/n)(1 - \ell/n)^{1/2 - \alpha} \leq 1 \) in general. Consequently condition (29) together with relation (30) finishes the proof. \( \square \)

For example, if \( \Delta \) is independent of \( n \), then condition (29) will hold for \( \ell \sim cn, c \in (0, 1) \). In case that \( n^\nu \ll \ell \ll n - n^\nu, \nu > 0 \), condition (29) holds provided that \( (1 - 2\alpha)/(1 - \alpha) < 2\nu \). That is, choosing \( \alpha \) close to \( 1/2 \) allows us to detect relatively “short” (“long”) epidemics.

6. Proof of the main theorems.

6.1. A moment inequality. In the proofs of our theorems we will use the following moment inequality which may be of separate interest.

**Proposition 4.** Let \( \{Y_k, k \in \mathbb{Z}\} \) be a centered stationary sequence, weakly \( \mathcal{M} \)-dependent in \( L^p \) with \( p \geq 2 \) and a rate function \( \delta(\cdot) \) satisfying

\[
D_p := \sum_{m=0}^{\infty} \delta(m) < \infty.
\]

Then for any \( n \in \mathbb{N}, b \in \mathbb{Z} \) we have

\[
E \left| \sum_{k=b+1}^{b+n} Y_k \right|^p \leq C_p n^{p/2},
\]

where \( C_p \) is a constant depending on \( p \) and the sequence \( \{Y_k\} \).

**Proof.** By stationarity, we can assume \( b = 0 \). Let first \( p = 2 \). We use below that \( \sup_{m \geq 0} \|Y_{k}^{(m)}\|_p \leq \|Y_1\|_p + D_p \). Without loss of generality we assume that \( EY_k^{(m)} = 0 \) for all \( k \in \mathbb{Z} \) and \( m \in \mathbb{N} \). Since

\[
Y_kY_{k+j} = (Y_k - Y_k^{(j-1)})Y_{k+j} + Y_k^{(j-1)}(Y_{k+j} - Y_{k+j}^{(j-1)})
\]

\[
+ Y_k^{(j-1)}Y_{k+j}^{(j-1)},
\]

we get by assumption (B) that for \( j \geq 1 \)

\[
|EY_kY_{k+j}| \leq |E[(Y_k - Y_k^{(j-1)})Y_{k+j}]| + |E[Y_k^{(j-1)}(Y_{k+j} - Y_{k+j}^{(j-1)})]|
\]

\[
\leq \|Y_{k+j}\|_2\|Y_k - Y_k^{(j-1)}\|_2 + \|Y_k^{(j-1)}\|_2\|Y_{k+j} - Y_{k+j}^{(j-1)}\|_2
\]

\[
\leq (\|Y_{k+j}\|_2 + \|Y_k^{(j-1)}\|_2)\delta(j-1)
\]

\[
\leq (2\|Y_1\|_2 + D_2)\delta(j-1).
\]
From relation (32) we infer, letting \( S_n = \sum_{k=1}^{n} Y_k \),

\[
ES_n^2 = \sum_{k=1}^{n} EY_k^2 + 2 \sum_{1 \leq k < l \leq n} EY_k Y_l \\
\leq n\|Y_1\|^2 + 2 \left[ \sum_{1 \leq k \leq n-1} |EY_k Y_{k+1}| + \cdots + \sum_{1 \leq k \leq 2} |EY_k Y_{k+n-2}| + E|Y_1 Y_n| \right] \\
\leq n\|Y_1\|^2 + 2(2\|Y_1\|^2 + D_2)[(n-1)\delta(0) + \cdots + 2\delta(n-3) + \delta(n-2)] \\
\leq n(\|Y_1\|^2 + 2D_2(2\|Y_1\|^2 + D_2)) =: C_2 n.
\]

This shows (31) for \( p = 2 \).

Once (31) is established for \( p \), it holds for all \( 0 < q \leq p \). Indeed, by Lyapunov’s inequality, relation (31) implies

\[
E \left| \sum_{k=b+1}^{b+n} Y_k \right|^q \leq C_p^{q/p} n^{q/2}
\]

for any \( 0 < q \leq p \). In particular, (31) holds with \( p = 1 \).

Next we prove (31) for all integers \( p > 2 \). Clearly, if \( C_p \geq \|Y_1\|^p \), then the inequality

\[
E|S_n|^p \leq C_p n^{p/2}
\]

holds for \( n = 1 \). Using a double induction argument, we show now that for some constant \( C_p \), relation (34) holds for all \( n \in \mathbb{N} \). More precisely, we show that if (34) holds for \( p-1 \) and all \( n \in \mathbb{N} \) and also for \( p \) and \( n \leq n_0 \), then it will also hold for \( p \) and \( n \leq 2n_0 \).

For \( k \leq n \) put \( S^n_k = Y_k + Y_{k+1} + \cdots + Y_n \). We have

\[
E|S_{2n}|^p = E|S_n + S_{2n+1}|^p \\
= E \left( \sum_{k=1}^{n} (Y_k - Y_{k+n-k}^{(n-k)}) + \sum_{k=1}^{n} (Y_{n+k} - Y_{n+k}^{(k-1)}) \right) \\
\quad + \sum_{k=1}^{n} Y_k^{(n-k)} + \sum_{k=1}^{n} Y_{n+k}^{(k-1)} \\
\leq \left( \sum_{k=1}^{n} \|Y_k - Y_{k}^{(n-k)}\|^p + \sum_{k=1}^{n} \|Y_{n+k} - Y_{n+k}^{(k-1)}\|^p \right) \\
\quad + \left( \sum_{k=1}^{n} Y_k^{(n-k)} + \sum_{k=1}^{n} Y_{n+k}^{(k-1)} \right)^p
\]

(35)
Due to assumption (B) in Definition

Hence (31)

By the induction assumption, this relation holds with arbitrary $n$ for all integer moments of order $≤ p$. The same estimate applies for $E|W_n|^p$. Due to assumption (B) in Definition 1, the random variables $Z_n$ and $W_n$ are independent. Thus

$$E|Z_n + W_n|^p$$

$$≤ n^{p/2} [2(1 + \psi_p)^p C_p + R_p].$$

Hence (36) and (37) and our assumptions on $C_p$ imply that

$$E|S_{2n}|^p ≤ (2\psi_p C_p^{1/p} + n^{1/2}[2(1 + \psi_p)^p C_p + R_p]^{1/p})^p$$

Choosing $\psi_p$ small enough, and then choosing $C_p$ large enough, we can always achieve that the term in brackets of (38) is $≤ \sqrt{2}$, provided that $p > 2$, and that the inequality $C_p^{1/p} > D_p/\psi_p$ mentioned before is satisfied. Hence we have for every $n ≤ n_0$ that $E|S_{2n}|^p ≤ C_p(2n)^{p/2}$, proving (34) for all even numbers $n ≤ 2n_0$. The case of odd $n$ is similar. The proof of Proposition 4 is finished for integer $p$.

For general $p > 2$ we have by the result shown before that (31) holds for $[p]$. (As usual, $[p]$ denotes the integer part of the real number $p$.) To finish the proof we need the following inequality which will be proven below:

$$|a + b|^p ≤ |a|^p + |b|^p$$

$$+ \sum_{k=1}^{[p]} \binom{p}{k} (|a|^k |b|^{p-k} + |b|^k |a|^{p-k}), \quad p ∈ [1, \infty).$$
Using (39) we get a similar estimate for $E|Z_n + W_n|^p$ as in (37) and the proof can be finished along the same lines as for integer $p$.

Verification of (39): Let $x \in [0, 1]$. We recall that $(1 + x)^p$ can be expanded in the binomial series

$$(1 + x)^p = \sum_{k \geq 0} \binom{p}{k} x^k$$

with

$$\binom{p}{k} = \frac{p(p-1) \cdots (p-k+1)}{k!}.$$  

From (40) it is clear that for $k \geq \lfloor p \rfloor + 2$ we have $\text{sign} \{ \binom{p}{k} \} = (-1)^{k-\lfloor p \rfloor + 1}$.

This immediately yields for $k = \lfloor p \rfloor + 2 \ell$ with $\ell \geq 1$,

$$\binom{p}{k} x^k + \binom{p}{k+1} x^{k+1} \leq \binom{p}{k} x^k + \binom{p}{k+1} x^k = \binom{p+1}{k+1} x^k < 0.$$  

Consequently

$$\sum_{k \geq \lfloor p \rfloor + 2} \binom{p}{k} x^k < 0$$

and

$$(1 + x)^p \leq \sum_{k=0}^{\lfloor |p| \rfloor + 1} \binom{p}{k} x^k.$$  

Now consider $|a + b|^p$. If $|a| \geq |b|$, then we infer from (41) that

$$|a + b|^p \leq |a|^p \left(1 + \frac{|b|}{|a|}\right) \leq |a|^p \sum_{k=0}^{\lfloor |p| \rfloor + 1} \binom{p}{k} \left(\frac{|b|}{|a|}\right)^k$$

$$= |a|^p + \sum_{k=1}^{\lfloor |p| \rfloor} \binom{p}{k} |b|^k |a|^{p-k}$$

$$+ \binom{p}{\lfloor |p| \rfloor + 1} |b|^p \left|\frac{b}{a}\right|^{\lfloor |p| \rfloor + 1 - p}.$$  

Thus (39) follows from $\binom{p}{\lfloor |p| \rfloor + 1} |b|^p |a|^{\lfloor |p| \rfloor + 1 - p} \leq 1$. Interchanging the roles of $a$ and $b$ completes the proof. \(\square\)

Using Möricz [33], Theorem 1, we get:
Corollary 1. Under the assumptions of Proposition 4 with \( p > 2 \), we have for any \( 2 < q \leq p \) and any \( n \in \mathbb{N}, b \in \mathbb{Z} \)

\[
E \max_{1 \leq k \leq n} \left| \sum_{j=b+1}^{b+k} Y_j \right|^q \leq C'_{p,q} n^{q/2},
\]

where the constants \( C'_{p,q} \) only depend on \( p,q \) and the sequence \( \{Y_k\} \).

A slightly weaker result can also be derived from Proposition 4 for the case of \( 0 < q \leq 2 \).

6.2. Proofs of Theorems 1 and 2. We give the proof of Theorem 1. Note first of all that \( \delta(m) = \|Y_k - Y_k^{(m)}\|_p \geq \|Y_k - Y_k^{(m)}\|_2 \), and consequently (32) holds when the \( L^2 \)-norm is replaced by the \( L^p \)-norm. Since \( A > 1 \) in (9), we infer that the series in (10) is absolutely convergent.

Let us specify some constants that will be used for the proof. By our assumption on \( A \) it is possible to find a constant \( 0 < \varepsilon_0 < 1/2 \) such that

\[
A > \frac{p - 2}{2\eta(1 - \varepsilon_0)} \left( 1 - \frac{1 + \eta}{p} \right).
\]

Then we set

\[
\delta = \frac{\beta}{1 + \alpha} \quad \text{with} \quad \alpha = \frac{2\eta(1 - \varepsilon_0)}{p - 2(1 + \eta)},
\]

(42)

\[
\beta = (1 - \varepsilon_0)\alpha.
\]

For some \( \varepsilon_1 > 0 \) (which will be specified later) we now define \( m_k = \lfloor \varepsilon_1 k^\delta \rfloor \).

The first step in the proof of (11) is to show that it is sufficient to provide the strong approximation for the perturbed sequence \( Y_k' = Y_k^{(m_k)} \). We notice that our main assumption (8) yields \( \|Y_k - Y_k'\|_p \ll k^{-A\delta} \). If \( A\delta < 1 \), then

\[
P \left( \max_{2^n \leq k \leq 2^{n+1}} \left| \sum_{j=1}^{k} (Y_j - Y_j') \right| > \frac{1}{n} 2^{(n/p)(1+\eta)} \right)
\]

\[
\leq P \left( \sum_{j=1}^{2^{n+1}} |Y_j - Y_j'| > \frac{1}{n} 2^{(n/p)(1+\eta)} \right)
\]

\[
\leq 2^{-(n(1+\eta))} n^p \left( \sum_{j=1}^{2^{n+1}} \|Y_j - Y_j'\|_p \right)^p
\]

\[
\ll 2^{-c_1 n} n^p,
\]
where \( c_1 = (1 + \eta) - (1 - A\delta)p > 0 \). Thus by the Borel–Cantelli lemma we have almost surely
\[
\sum_{j=1}^{k} Y_j = \sum_{j=1}^{k} Y'_j + o(k^{(1+\eta)/p}) \quad \text{a.s.}
\]
If \( A\delta \geq 1 \) we get an (even better) error term of order \( o(k^{1/p}) \).

The main part of the proof of Theorem 1 is based on a blocking argument. We partition \( \mathbb{N} \) into disjoint blocks
\[
\mathbb{N} = J_1 \cup I_1 \cup J_2 \cup I_2 \cup \cdots,
\]
where \( |I_k| = \lfloor k^\alpha \rfloor \) and \( |J_k| = \lfloor k^\beta \rfloor \) with \( \alpha, \beta \) as in (42). Let us further set
\[
I_k = \{ i_k, \ldots, \bar{i}_k \} \quad \text{and} \quad J_k = \{ j_k, \ldots, \bar{j}_k \}
\]
and
\[
\xi_k = \sum_{j \in I_k} Y'_j \quad \text{and} \quad \eta_k = \sum_{j \in J_k} Y'_j.
\]
Note that \( \bar{i}_k = O(k^{1+\alpha}) \). Provided that \( \varepsilon_1 \) in the definition of \( m_k \) is chosen small enough, this will imply that
\[
|J_k| = \lfloor k^\beta \rfloor > \lfloor \varepsilon_1 k^\delta \rfloor = m_{\bar{k}},
\]
and hence by assumption (B) it follows that \( \{\xi_k\} \) and \( \{\eta_k\} \) each define a sequence of independent random variables.

The following lemma by Sakhanenko [43] (cf. also Shao [44]) is our crucial ingredient for the construction of the approximating processes.

**Lemma 2.** Let \( \{\xi_k\} \) be a sequence of centered independent random variables with finite \( p \)th moments, \( p > 2 \). Then we can redefine \( \{\xi_k\} \) on a suitable probability space, together with a sequence \( \{\xi_k^*\} \) of independent normal random variables with \( E\xi_k^* = 0 \), \( E(\xi_k^*)^2 = E\xi_k^2 \) such that for any \( x > 0 \), \( m \geq 1 \)
\[
P\left( \max_{1 \leq k \leq m} \left| \sum_{j=1}^{k} \xi_j - \sum_{j=1}^{k} \xi_j^* \right| > x \right) \leq C \frac{1}{x^p} \sum_{j=1}^{m} E|\xi_j|^p,
\]
where \( C \) is an absolute constant.

We shall now apply Lemma 2 to the sequences \( \{\xi_k\} \) and \( \{\eta_k\} \). For this purpose we need estimates of the moments \( E|\xi_k|^p, E|\eta_k|^p \). By Minkowski’s inequality and Proposition 4 we get
\[
E|\xi_k|^p \leq \left( \left\| \sum_{j \in I_k} Y_j \right\|_p + \sum_{j \in I_k} \left\| Y_j - Y'_j \right\|_p \right)^p
\]
\[
= O\left( (|I_k|^{1/2} + |I_k| \cdot \bar{i}_k^{-A\delta})^p \right).
\]
Some easy algebra shows that the restrictions on the parameters $A$, $\delta$, $\alpha$ and $\varepsilon_0$ imply

$$|I_k| \cdot k^{-A\delta} \ll k^\alpha \cdot k^{-Al(1+\alpha)} \ll k^{\alpha / 2} \ll |I_k|^{1 / 2}.$$  

A similar estimate holds for $E|\eta_k|^p$. Hence we can find constants $F_p$ such that

$$E|\xi_k|^p \leq F_p |I_k|^{p / 2}$$

and

$$E|\eta_k|^p \leq F_p |J_k|^{p / 2},$$

where $F_p$ does not depend on $k$.

Let $L_n = \sum_{k=1}^n |I_k|$. Then $L_n = O(n(1+\alpha))$. By our previous estimates and by Lemma 2 we infer that, after enlarging the probability space, we have

$$P \left( \max_{2^n \leq k \leq 2^{n+1}} \left| \sum_{j=1}^k \xi_j - \sum_{j=1}^k \xi_j^* \right| > L_{2^n}^{(1+\eta)/p} \right)$$

$$\leq L_{2^n}^{-(1+\eta)} \sum_{k=1}^{2^{n+1}} E|\xi_k|^p$$

$$= O(2^{-(1+\alpha)(1+\eta)+\alpha p/2+1}) n),$$

where $\xi_k^*$ is a sequence of independent and centered normal random variables with $E(\xi_k^*)^2 = E\xi_k^2$. The exponent in (43) will be negative if $(1+\alpha)(1+\eta) > \frac{\alpha p}{2} + 1$. This is equivalent to $\alpha < \frac{2^{p}(1+\eta)}{2^{p}(1+\eta)}$, which follows by (42). Thus, by the Borel–Cantelli lemma we obtain

$$\sum_{j=1}^k \xi_j = \sum_{j=1}^k \xi_j^* + O(L_k^{(1+\eta)/p}) \quad \text{a.s.}$$

By further enlarging the probability space we can write

$$\sum_{j=1}^k \xi_j = W_1 \left( \sum_{j=1}^k \operatorname{Var}(\xi_j) \right) + O(L_k^{(1+\eta)/p}) \quad \text{a.s.,}$$

where $\{W_1(t), t \geq 0\}$ is a standard Wiener process. The same arguments show that

$$\sum_{j=1}^k \eta_j = W_2 \left( \sum_{j=1}^k \operatorname{Var}(\eta_j) \right) + O(M_k^{(1+\eta)/p}) \quad \text{a.s.,}$$

where $\{W_2(t), t \geq 0\}$ is another standard Wiener process on the same probability space and $M_n = \sum_{k=1}^n |J_k|$.
We define
\[ b^2_k = \frac{\text{Var}\left( \sum_{j \in I_k} Y'_j \right)}{|I_k|} \]
and
\[ h^2_k = \frac{\text{Var}\left( \sum_{j \in J_k} Y'_j \right)}{|J_k|}. \]

For \( \ell \in I_k \) we set \( \sigma^2_{\ell} = b^2_k \) and for \( \ell \in J_k \) we set \( \sigma^2_{\ell} = 0. \) Similarly define \( \tau^2_{\ell} = h^2_k \) if \( \ell \in J_k \) and \( \tau^2_{\ell} = 0 \) if \( \ell \in I_k. \) Put
\[ s^2_n = \sum_{k=1}^{n} \sigma^2_k, \quad t^2_n = \sum_{k=1}^{n} \tau^2_k \quad (n = 1, 2, \ldots). \]

Summarizing our results so far we can write
\[ \sum_{k=1}^{n} Y_k = W_1 \left( \sum_{k=1}^{n} \sigma_k^2 \right) + W_2 \left( \sum_{k=1}^{n} \tau_k^2 \right) + O\left( \frac{1}{n^{1+\eta}} \right) \quad \text{a.s.} \]

It is a basic result that our stationarity and dependence assumptions imply
\[
\text{Var}\left( \sum_{j \in I_k} Y_j \right)/|I_k| = \sigma^2 + O(k^{-\xi}) \quad \text{and}
\]
\[
\text{Var}\left( \sum_{j \in J_k} Y_j \right)/|J_k| = \sigma^2 + O(k^{-\xi}) \tag{44}
\]
as \( k \to \infty, \) for some small enough \( \xi > 0. \) It can be easily shown that (44) remains true if the \( Y_j \) are replaced with \( Y'_j. \) Indeed, by the Minkowski inequality we infer that
\[
\text{Var}^{1/2} \left( \sum_{j \in I_k} Y'_j \right) \leq \text{Var}^{1/2} \left( \sum_{j \in I_k} Y'_j \right) + \text{Var}^{1/2} \left( \sum_{j \in I_k} (Y_j - Y'_j) \right)
\]
\[
\leq \text{Var}^{1/2} \left( \sum_{j \in I_k} Y'_j \right) + |I_k| \max_{j \in I_k} \|Y_j - Y'_j\|_2.
\]

Furthermore, using the definitions of the introduced constants we obtain
\[
\max_{j \in I_k} \|Y_j - Y'_j\|_2 \ll \|J_k\|^{-A\delta} \ll k^{-(\alpha+1)A\delta}
\]
\[
\leq k^{-\beta} = k^{-(1-\epsilon_0)\alpha} \quad \text{with } \epsilon_0 < 1/2.
\]
Since by definition \( |I_k| \ll k^\alpha \), we conclude that
\[
\text{Var}^{1/2} \left( \sum_{j \in I_k} Y_j' \right) / |I_k|^{1/2} \leq \text{Var}^{1/2} \left( \sum_{j \in I_k} Y_j \right) / |I_k|^{1/2} + O\left( k^{\alpha(\epsilon_0 - 1/2)} \right)
\]
as \( k \to \infty \). In the same manner a lower bound for \( \text{Var}^{1/2} \left( \sum_{j \in I_k} Y_j' \right) / |I_k|^{1/2} \) can be obtained. Proving the analogue of the second part of (44) for the \( Y_j' \) is similar.

In other words, we have shown (11) along the subsequence \( \{ i_n \} \) with values of \( s^2_n \) and \( t^2_m \) that satisfy (12) and (13). The relation \( |s^2_n - \sigma^2 n| = O(n^{1-\epsilon}) \), \( \epsilon > 0 \), follows by simple calculations.

To finish the proof we have to show that the fluctuations of the partial sums and the Wiener processes \( W_1 \) and \( W_2 \) within the blocks \( I_k \) are small enough. Since fluctuation properties of Wiener processes are easy to handle using standard deviation inequalities (see, e.g., [9]), we only investigate the partial sums. By Corollary 1 we have
\[
P \left( \sup_{i_k \leq \ell \leq i_k} \left| \sum_{j = i_k}^\ell Y_j \right| > \frac{\text{L}}{k}^{(1+\eta)/p} \right) \leq \frac{\text{L}}{k}^{(1+\eta)} E \left( \sup_{i_k \leq \ell \leq i_k} \left| \sum_{j = i_k}^\ell Y_j \right|^p \right)
\]
\[
\ll \frac{\text{L}}{k}^{(1+\eta)} |I_k|^{p/2}
\]
\[
\ll k^{-(1+\eta)(1+\alpha)+ap/2}
\]
\[
= O(k^{-(1+\epsilon_2)}),
\]
if \( \epsilon_2 > 0 \) is chosen sufficiently small. The Borel–Cantelli lemma shows that we can also control the fluctuation within the blocks. Thus (11) is proven.

The proof of Theorem 2 is similar to the proof of Theorem 1 and will be therefore omitted. We only remark that under the exponential mixing rate logarithmic block sizes are required in the blocking argument.

6.3. **Proof of Proposition 1.** We use the notation introduced in the proof of Theorem 1. Further we let \( I = I_1 \cup I_2 \cup \cdots \) and \( J = J_1 \cup J_2 \cup \cdots \) and \( M_n = \{1, \ldots, n\} \). By looking at the proof of Theorem 1, it readily follows that
\[
\frac{1}{s_n} \sum_{i \in I \cap M_n} Y_i = W_1(s^2_n)/s_n - X_n,
\]
\[
\frac{1}{t_m} \sum_{j \in J \cap M_n} Y_j = W_2(t^2_m)/t_m - Z_m,
\]
where
\[
X_n = o\left( (s^2_n)^{(1+\eta)/p-1/2} \right) = o(1) \quad \text{a.s. and}
\]
\[
Z_m = o\left( (t^2_m)^{(1+\eta)/p-1/2} \right) = o(1) \quad \text{a.s.}
\]
(45)
Hence
\[
\text{Corr}(W_1(s_n^2), W_2(t_m^2)) = \text{Corr}\left(\frac{1}{s_n} W_1(s_n^2), \frac{1}{t_m} W_2(t_m^2)\right)
= \text{Corr}\left(\frac{1}{s_n} \sum_{i \in I \cap M_n} Y_i + X_n, \frac{1}{t_m} \sum_{j \in J \cap M_m} Y_j + Z_m\right).
\]

In order to calculate this correlation we need a couple of estimates.

First we note that by the definition of \(s_n^2\) and \(t_m^2\)
\[
s_n^2 \sim \sigma^2 |I \cap M_n| \quad \text{and} \quad t_m^2 \sim \sigma^2 |J \cap M_m|.
\]

It readily follows from Proposition 4 that
\[
\left\| \frac{1}{s_n} \sum_{i \in I \cap M_n} Y_i \right\|_p \leq C_p,
\]
where \(C_p\) does not depend on \(n\). Thus
\[
\sup_{n \geq 1} \left\| X_n \right\|_p = \sup_{n \geq 1} \left\| \frac{1}{s_n} \sum_{i \in I \cap M_n} Y_i - W_1(s_n^2)/s_n \right\|_p
\leq \sup_{n \geq 1} \left\| \frac{1}{s_n} \sum_{i \in I \cap M_n} Y_i \right\|_p + \left\| W_1(1) \right\|_p < \infty,
\]
and hence \(\{X_n^2\}\) is uniformly integrable. This and (45) show that \(\text{Var}(X_n) \to 0\); by the same arguments \(\text{Var}(Z_m) \to 0\). By (44)
\[
\left\| \frac{1}{s_n} \sum_{i \in I \cap M_n} Y_i \right\|_2 \sim \sigma^2.
\]
Thus by (47) and (48)
\[
c_1(m, n) := \text{Cov}\left( \frac{1}{s_n} \sum_{i \in I \cap M_n} Y_i \right)
\leq \text{Var}^{1/2}(Z_m) \text{Var}^{1/2}\left( \frac{1}{s_n} \sum_{i \in I \cap M_n} Y_i \right)
= o(1) \quad \text{for } m, n \to \infty,
\]
and similarly
\[
c_2(m, n) := \text{Cov}\left( \frac{1}{t_m} \sum_{j \in J \cap M_m} Y_j \right) = o(1) \quad \text{for } m, n \to \infty.
\]
Furthermore we have

\[ B_1(n) := \text{Var}^{1/2}\left(\frac{1}{s_n} \sum_{i \in I \cap M_n} Y_i + X_n\right) \]

\[ \geq \text{Var}^{1/2}\left(\frac{1}{s_n} \sum_{i \in I \cap M_n} Y_i\right) - \text{Var}^{1/2}(X_n) \]

\[ = \sigma + o(1) \quad \text{for } n \to \infty \]

and

\[ B_2(m) := \text{Var}^{1/2}\left(\frac{1}{t_m} \sum_{j \in J \cap M_m} Y_j + Z_m\right) \]

\[ \geq \sigma + o(1) \quad \text{for } m \to \infty. \]

Finally we introduce the term

\[ c_0(m, n) = \frac{1}{s_n t_m} \sum_{i \in I \cap M_n} \sum_{j \in J \cap M_m} \text{Cov}(Y_i, Y_j). \]

We choose \( r \geq 0 \) such that \( n \in I_{r+1} \cup J_{r+1} \), and we choose \( v \geq 0 \) such that \( m \in I_{v+1} \cup J_{v+1} \) and recall that by Theorem 1 we have \( \sum_{i \in \mathbb{Z}} |\text{Cov}(Y_0, Y_i)| < \infty \). Hence if \( v \leq 2r \) we have

\[ c_0(m, n) \leq s_{tr}^{-1} t_{rv}^{-1} \sum_{i \in I_{1} \cup \cdots \cup I_{r+1}} \sum_{j \in J_{1} \cup \cdots \cup J_{v+1}} |\text{Cov}(Y_i, Y_j)| \]

\[ \leq s_{tr}^{-1} t_{rv}^{-1} \sum_{j \in J_{1} \cup \cdots \cup J_{v+1}} \sum_{i \in \mathbb{Z}} |\text{Cov}(Y_i, Y_{0})| \]

\[ = s_{tr}^{-1} t_{rv}^{-1} \sum_{j \in J_{1} \cup \cdots \cup J_{v+1}} \sum_{i \in \mathbb{Z}} |\text{Cov}(Y_i, Y_{0})| \]

\[ \ll s_{tr}^{-1} t_{rv}^{-1} (|J_{1}| + \cdots + |J_{v+1}|) \]

\[ \ll s_{tr}^{-1} t_{rv}^{-1} = o(1) \quad \text{as } m, n \to \infty. \]

If \( v > 2r \), we have to additionally show that

\[ s_{tr}^{-1} t_{rv}^{-1} \sum_{i \in I_{1} \cup \cdots \cup I_{r+1}} \sum_{j \in J_{2r+1} \cup \cdots \cup J_{v+1}} |\text{Cov}(Y_i, Y_j)| \rightarrow 0. \]

Now we have by (32) and assumptions (8), (9) that

\[ s_{tr}^{-1} t_{rv}^{-1} \sum_{i \in I_{1} \cup \cdots \cup I_{r+1}} \sum_{j \in J_{2r+1} \cup \cdots \cup J_{v+1}} |\text{Cov}(Y_i, Y_j)| \]

\[ \ll s_{tr}^{-1} t_{rv}^{-1} = o(1) \quad \text{as } m, n \to \infty. \]
\[ \leq s^{-1}_{ir}t^{-1}_{iv} \sum_{\pi \geq 2r+1} |J_{\pi}| \sum_{\ell=1}^{r+1} |I_{\ell}| (d(I_{\ell}, J_{\pi}))^{-1} \]

\[ \ll s^{-1}_{ir}t^{-1}_{iv} \sum_{\pi \geq 2r+1} \pi^\beta \sum_{\ell=1}^{r+1} \ell^\alpha (d(I_{\ell}, J_{\pi}))^{-1}. \]

For \( \ell \in \{1, \ldots, r+1\} \) and \( \pi \geq 2r+1 \) we have constants \( k_0 \) and \( k_1 \) independent of \( r \) and \( \pi \) such that
\[
d(I_{\ell}, J_{\pi}) \geq k_0 (\pi^\alpha - r^\alpha) \geq k_1 \pi^\alpha+1
\]
and thus
\[
s^{-1}_{ir}t^{-1}_{iv} \sum_{\pi \geq 2r+1} \pi^\beta \sum_{\ell=1}^{r+1} \ell^\alpha (d(I_{\ell}, J_{\pi}))^{-1} \ll s^{-1}_{ir}t^{-1}_{iv} \pi^\alpha+1 \sum_{\pi \geq 2r+1} \pi^{\beta-\alpha-1} \ll s^{-1}_{ir}t^{-1}_{iv} r^{1+\beta} \ll r^{-\alpha-\beta/2} = o(1) \quad \text{as } r \to \infty.
\]

Using the definitions of \( c_0, c_1, c_2 \) and \( B_1 \) and \( B_2 \) we see that
\[
\text{Corr}(W_1(s^2_n), W_2(t^2_m)) = \frac{c_0(m, n) + c_1(m, n) + c_2(m, n) + \text{cov}(X_n, Z_m)}{B_1(n)B_2(m)}.
\]

We have shown \( c_0(m, n) + c_1(m, n) + c_2(m, n) + \text{cov}(X_n, Z_m) \to 0 \) as \( m, n \to \infty \) while the denominator in (49) is bounded away from zero. This finishes the proof of Proposition 1.

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