# Quadratic covariation estimation of an irregularly observed semimartingale with jumps and noise 

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#### Abstract

This paper presents a central limit theorem for a pre-averaged version of the realized covariance estimator for the quadratic covariation of a discretely observed semimartingale with noise. The semimartingale possibly has jumps, while the observation times show irregularity, non-synchronicity, and some dependence on the observed process. It is shown that the observation times' effect on the asymptotic distribution of the estimator is only through two characteristics: the observation frequency and the covariance structure of the noise. This is completely different from the case of the realized covariance in a pure semimartingale setting.


Keywords: jumps; microstructure noise; non-synchronous observations; quadratic covariation; stable limit theorem; time endogeneity

## 1. Introduction

The quadratic covariation matrix of a semimartingale is one of the fundamental quantities in statistics of semimartingales. In the context of the estimation of the diffusion coefficient of an Itô process observed discretely in a fixed interval, limit theorems associated with the discretized quadratic covariation play a key role, and such research has a long history (cf. [17,19]). Furthermore, in recent years such an asymptotic theory has been applied to measuring the covariance structure of financial assets from high-frequency data. This was pioneered by [4,7], and has become one of the most active areas in financial econometrics. In such a context, the discretized quadratic covariation is also called the realized covariance.

However, raw high frequency data typically deviates from the ideal situation where we observe a continuous semimartingale at equidistant times, and this motivates statisticians to develop the theory in more complicated settings. One topic is the treatment of measurement errors in the data. For financial high-frequency data, such errors originate from market microstructure noise and have attracted vast attention in the past decade; among various studies, see, for example, [5,46,47,52,54]. In the univariate context, central limit theorems under irregular sampling settings have also been studied by many authors, especially assuming the independence between the observed process and the observation times; see, for example, [19,22,42]. In the multivariate case, the irregularity of the observation times causes the non-synchronicity which makes the analysis more complicated. The prominent works on this topic are the Fourier analysis approach of [40], the sampling design kernel method of [23] and the quasi-likelihood analysis of [44]. In addition,
recently various approaches to deal with these issues simultaneously have been proposed by many authors; see, for example, [1,6,8,9,15,51].

Another important issue is incorporating jumps into the model. In such a situation interest is often paid to estimating the integrated volatility and the integrated covariance matrix, that is, the integrated diffusion coefficient, and there are many studies on this issue in various settings. Regarding the central limit theories, see, for example, Chapters 11 and 13 of [31] for the basic setting, [45] for the noise setting, [34] for the non-synchronous observation setting and [13] for the noisy and non-synchronous observation setting.

In contrast, turning to the entire quadratic covariation estimation in the presence of jumps, there are fewer works. A central limit theorem for the realized covariance of an equidistantly observed Lévy process has been proved in [30] in the context of the analysis of the Euler scheme. This result has been extended to general Itô semimartingales in [27] as a special case of the asymptotic results on various functionals of semimartingale increments. The situation where measurement errors are present has been treated by Jacod, Podolskij and Vetter [29] who focus on the "pre-averaging" counterparts of the functionals discussed in [27], which were introduced in [46] to extend classical power variation based methods to a noisy observation setting. The theory requires a different treatment in the absence of the diffusion coefficient, and this case has been studied in [16].

When we further focus on the situation where the observation times are irregular, at least to the best of the author's knowledge, there is no comprehensive study on the central limit theory for the quadratic covariation estimation, except for the recent work of Bibinger and Vetter [12] and Bibinger and Winkelmann [13]; the former have derived central limit theorems for the realized covariance and the Hayashi-Yoshida estimator of Hayashi and Yoshida [23] for a general Itô semimartingale observed irregularly and non-synchronously, while the latter have established a central limit theorem for an adjusted version of the spectral covariance estimator of Bibinger and Reiß [11] in a non-synchronous and noisy observation setup, focusing on asymptotically regular observation times in the sense that they satisfy conditions in Proposition 2.54 of [43]. The aim of this study is to develop such a theory in the situation where the observation data is contaminated by noise and the observation times are as general as possible. More precisely, we derive a central limit theorem for the pre-averaged version of the realized covariance proposed in [14] (called the modulated realized covariance) under an irregular sampling setting in the presence of jumps. The main finding of this paper is that in the synchronous case the observation times' effect on the asymptotic distribution of the estimator is only through their conditional expected durations, provided that the limit of such quantities are well-defined. In other words, the irregularity of the observation times has no impact on the asymptotic distribution of the estimator because the conditional expected durations of the observation times naturally link with the magnitude of the observation frequency, and thus their effect is not due to the irregularity. This is completely different from the pure semimartingale setting of [12] where the distribution of the durations around the jump times of the semimartingale directly affects the asymptotic distribution of the realized covariance.

To deal with non-synchronous observations we rely on a data synchronization method proposed in [1], which also matches the proposal of Section 3.6 of [14]. The non-synchronicity naturally links with the covariance structure of the noise, hence it affects the asymptotic distribution through that relation. On the other hand, the interpolations to the synchronized sampling
times do not matter asymptotically. This can be seen as a counterpart of the finding of [8] in the continuous case.

Another issue we attempt to solve is how the dependence between the observed process and the observation times (called the time endogeneity) affects the asymptotic theory in our setting. This issue has recently been highlighted by several authors such as Fukasawa [18], Li et al. [37], Li, Zhang and Zheng [38], Rosenbaum and Tankov [50] in various settings, and it is indeed known that such dependence possibly causes a non-standard limit theorem even in the continuous semimartingale setting. In this paper, this issue is partly solved in the sense that we do not rule out the dependence between the continuous component of the process and the observation times, but partly rule out the dependence between the jump component and the observation times. The result shows that the time endogeneity is also immaterial in our setting.

This paper is organized as follows. Section 2 presents the mathematical model and the construction of the estimator we are focusing on. Section 3 is devoted to the main result of this paper. Section 4 provides some illustrative examples of the observation times, while Section 5 provides a simulation study. All proofs are given in Section 6.

## 2. The set up

Given a stochastic basis $\mathcal{B}^{(0)}=\left(\Omega^{(0)}, \mathcal{F}^{(0)},\left(\mathcal{F}_{t}^{(0)}\right)_{t \geq 0}, P^{(0)}\right)$, we consider a $d$-dimensional semimartingale $X=\left(X_{t}\right)_{t \in \mathbb{R}_{+}}$of the form

$$
X_{t}=X_{0}+\int_{0}^{t} b_{s} \mathrm{~d} s+\int_{0}^{t} \sigma_{s} \mathrm{~d} W_{s}+\left(\delta 1_{\{\|\delta\| \leq 1\}}\right) \star(\mu-v)_{t}+\left(\delta 1_{\{\|\delta\|>1\}}\right) \star \mu_{t}
$$

where $W$ is a $d^{\prime}$-dimensional $\left(\mathcal{F}_{t}^{(0)}\right)$-standard Brownian motion, $\mu$ is an $\left(\mathcal{F}_{t}^{(0)}\right)$-Poisson random measure on $\mathbb{R}_{+} \times E$ with $E$ being a Polish space, $v$ is the intensity measure of $\mu$ of the form $\nu(\mathrm{d} t, \mathrm{~d} z)=\mathrm{d} t \otimes \lambda(\mathrm{~d} z)$ with $\lambda$ being a $\sigma$-finite measure on $E, b$ is an $\left(\mathcal{F}_{t}^{(0)}\right)$-progressively measurable $\mathbb{R}^{d}$-valued process, $\sigma$ is an $\left(\mathcal{F}_{t}^{(0)}\right)$-progressively measurable $\mathbb{R}^{d} \otimes \mathbb{R}^{d^{\prime}}$-valued process, and $\delta$ is an $\left(\mathcal{F}_{t}^{(0)}\right)$-predictable $\mathbb{R}^{d}$-valued function on $\Omega^{(0)} \times \mathbb{R}_{+} \times E$. Also, $\star$ denotes the integral (either stochastic or ordinary) with respect to some (integer-valued) random measure. Here and below, we use standard concepts and notation in stochastic calculus, which are described in detail in, for example, Chapter 2 of [31]. Our aim is to estimate the quadratic covariation matrix process $[X, X]=\left(\left[X^{k}, X^{l}\right]\right)_{1 \leq k, l \leq d}$ of $X$ from noisy and discrete observation data of $X$.

The observed process $Y$ is subject to additional measurement errors as follows:

$$
Y_{t}=X_{t}+\epsilon_{t}
$$

The mathematical construction of the noise process $\epsilon$ is explained later. We observe the components of the $d$-dimensional process $Y=\left(Y^{1}, \ldots, Y^{d}\right)$ discretely and non-synchronously. For each $k=1, \ldots, d$ the observation times for $Y^{k}$ are denoted by $t_{0}^{k}, t_{1}^{k}, \ldots$, that is, the observation data $\left(Y_{t_{i}^{k}}^{k}\right)_{i \in \mathbb{Z}_{+}}$is available. We assume that $\left(t_{i}^{k}\right)_{i=0}^{\infty}$ is a sequence of $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping times which implicitly depend on a parameter $n \in \mathbb{N}$ representing the observation frequency and satisfy that
$t_{i}^{k} \uparrow \infty$ as $i \rightarrow \infty$ and $\sup _{i \geq 0}\left(t_{i}^{k} \wedge t-t_{i-1}^{k} \wedge t\right) \rightarrow^{p} 0$ as $n \rightarrow \infty$ for any $t \in \mathbb{R}_{+}$, with setting $t_{-1}^{k}=0$ for notational convenience (hereafter we will refer to such a sequence as a sampling scheme for short).

Now we introduce the precise definition of the noise process $\epsilon$. It is basically the same as the one from Chapter 16 of [31], but we need a slight modification to ensure the (joint) measurability of the process $\epsilon$, which is necessary for us to consider variables such as $\epsilon_{t_{i}^{k}}^{k}$. For any $t \in \mathbb{R}_{+}$there is a transition probability $Q_{t}\left(\omega^{(0)}, \mathrm{d} u\right)$ from $\left(\Omega^{(0)}, \mathcal{F}_{t}^{(0)}\right)$ into $\mathbb{R}^{d}$ satisfying $\int u Q_{t}\left(\omega^{(0)}, \mathrm{d} u\right)=0$ (this will correspond to the conditional distribution of the noise at the time $t$ given $\left.\mathcal{F}_{t}^{(0)}\right)$. Then, at each frequency $n \in \mathbb{N}$, the stochastic basis $\mathcal{B}=\left(\Omega, \mathcal{F},\left(\mathcal{F}_{t}\right)_{t \in \mathbb{R}_{+}}, P\right)$ supporting the observed process $Y$ is constructed in the following manner (for notational simplicity we subtract the index $n$ from $\mathcal{B}$ ): we endow the space $\Omega^{(1)}=\left(\mathbb{R}^{d}\right)^{\mathbb{N}}$ with the product Borel $\sigma$-field $\mathcal{F}^{(1)}$ and with the probability measure $Q\left(\omega^{(0)}, \mathrm{d} \omega^{(1)}\right)$ which is the product $\bigotimes_{i \in \mathbb{N}} Q_{\mathcal{T}_{i}^{n}\left(\omega^{(0)}\right)}\left(\omega^{(0)}, \cdot\right)$. Here, $\left(\mathcal{T}_{i}^{n}\right)_{i \geq 0}$ is the increasing reordering of total observation times $\left\{t_{i}^{k}: k=1, \ldots, d\right.$ and $\left.i \in \mathbb{Z}_{+}\right\}$. More formally, it is defined sequentially by $\mathcal{T}_{0}^{n}=\min _{k=1, \ldots, d} t_{0}^{k}$ and $\mathcal{T}_{i}^{n}=\min _{k=1, \ldots, d} \min \left\{t_{j}^{k}: t_{j}^{k}>\mathcal{T}_{i-1}^{n}\right\}$ for $i=1,2 \ldots$. Note that $\mathcal{T}_{i}^{n}$ is an $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping time since $\mathcal{T}_{i}^{n}=\min _{k=1, \ldots, d} \inf _{j \geq 1}\left(t_{i}^{k}\right)_{\left\{t_{j}^{k}>\mathcal{T}_{i-1}^{n}\right\}}$, where for an $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping time $\tau$ and a set $A \in \mathcal{F}_{\tau}^{(0)}$, we define $\tau_{A}$ by $\tau_{A}\left(\omega^{(0)}\right)=\tau\left(\omega^{(0)}\right)$ if $\omega^{(0)} \in A ; \tau_{A}\left(\omega^{(0)}\right)=\infty$ otherwise (see Claim 1.15 from Chapter 1 of [32]). Then we define the probability space $(\Omega, \mathcal{F}, P)$ by

$$
\begin{align*}
\Omega & =\Omega^{(0)} \times \Omega^{(1)}, \quad \mathcal{F}=\mathcal{F}^{(0)} \otimes \mathcal{F}^{(1)}, \\
P\left(\mathrm{~d} \omega^{(0)}, \mathrm{d} \omega^{(1)}\right) & =P^{(0)}\left(\mathrm{d} \omega^{(0)}\right) Q\left(\omega^{(0)}, \mathrm{d} \omega^{(1)}\right) \tag{2.1}
\end{align*}
$$

After that, the noise process $\epsilon=\left(\epsilon_{t}\right)_{t \geq 0}$ is defined on this probability space by $\epsilon_{t}=\epsilon_{\mathrm{N}_{n}(t)}^{0}$, where $\left(\epsilon_{i}^{0}\right)_{i \in \mathbb{N}}$ denotes the canonical process on $\left(\Omega^{(1)}, \mathcal{F}^{(1)}\right)$ and $N_{n}(t)=\sum_{i=0}^{\infty} 1_{\left\{\mathcal{T}_{i}^{n} \leq t\right\}}$. By construction, given $\mathcal{F}^{(0)},\left(\epsilon_{\mathcal{T}_{i}^{n}}\right)_{i \in \mathbb{Z}_{+}}$is (serially) independent and $\epsilon_{\mathcal{T}_{i}^{n}}$ obeys the law $Q_{\mathcal{T}_{i}^{n}\left(\omega^{0}\right)}\left(\omega^{(0)}, \cdot\right)$ for every $i$. Finally, the filtration $\left(\mathcal{F}_{t}\right)_{t \geq 0}$ is defined as the one generated by $\left(\mathcal{F}_{t}^{(0)}\right)_{t \geq 0}$ and $\left(\epsilon_{t}\right)_{t \geq 0}$.

Any variable or process defined on either $\Omega^{(0)}$ or $\Omega^{(1)}$ is considered in the usual way as a variable or a process on $\Omega$. Specifically, our noisy process $Y=\left(Y_{t}\right)_{t \geq 0}$ is the process defined as the sum of the latent process $X$ on $\Omega^{(0)}$ and the noise process $\epsilon$ on $\Omega$.

Remark 2.1. To ensure that the probability measure $P$ in (2.1) is well-defined, we further need the measurability of the map $\omega^{(0)} \mapsto Q\left(\omega^{(0)}, A\right)$ for any Borel subset $A$ of $\mathbb{R}^{d}$. This is ensured by the progressive measurability of the process $\left(Q_{t}(\cdot, A)\right)_{t \geq 0}$ which we will assume later (see assumption [A4]). This assumption also ensures that $\mathcal{B}$ is the very good filtered extension of $\mathcal{B}^{(0)}$, which is necessary to apply the version of Jacod's stable limit theorem described by Theorem 2.2.15 of [31].

To deal with the non-synchronicity of the observation times we rely on a data synchronization method, which is commonly used in the literature; see, for example, [1,6,14,53]. Let $\left(T_{p}\right)_{p=0}^{\infty}$
and $\left(\tau_{p}^{k}\right)_{p=0}^{\infty}(k=1, \ldots, d)$ be sampling schemes such that

$$
\begin{equation*}
\tau_{0}^{k} \leq T_{0} \quad \text { and } \quad T_{p-1}<\tau_{p}^{k} \leq T_{p} \quad \text { for any } p \geq 1 \text { and any } k=1, \ldots, d \tag{2.2}
\end{equation*}
$$

We assume that the observation data $\left(Y_{\tau_{p}^{k}}^{k}\right)_{p \in \mathbb{Z}_{+}}$is available for every $k=1, \ldots, d$, that is, $\left\{\tau_{p}^{k}: p \geq 0\right\} \subset\left\{t_{i}^{k}: i \geq 0\right\}$, and construct statistics based on this synchronized data set $\left(Y_{\tau_{p}^{k}}^{k}\right)_{p \in \mathbb{Z}_{+}}$, $k=1, \ldots, d$. In [1] this type of synchronization method is called the Generalized Synchronization method and $\left(T_{p}\right)_{p=0}^{\infty}$ is called the Generalized Sampling Time. One way to implement such synchronization is the so-called refresh time sampling method introduced by Barndorff-Nielsen et al. [6] to this area. Namely, we first define the refresh times $T_{0}, T_{1}, \ldots$ of the sampling schemes $\left\{\left(t_{i}^{k}\right)\right\}_{k=1}^{d}$ sequentially by $T_{0}=\max \left\{t_{0}^{1}, \ldots, t_{0}^{d}\right\}$ and $T_{p}=\max _{k=1, \ldots, d} \min \left\{t_{i}^{k}: t_{i}^{k}>T_{p-1}\right\}$ for $p=1,2, \ldots$. After that, for each $k,\left(\tau_{p}^{k}\right)$ is defined by interpolating the next-ticks into $\left(T_{p}\right)$ as follows:

$$
\tau_{0}^{k}=t_{0}^{k} \quad \text { and } \quad \tau_{p}^{k}=\min \left\{t_{i}^{k}: t_{i}^{k}>T_{p-1}\right\}, \quad p=1,2, \ldots
$$

Note that $\tau_{p}^{k}$ is an $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping time due to an analogous reason to that for $\mathcal{T}_{i}^{n}$.
Now the modulated realized covariance (henceforth MRC) estimator we focus on is constructed in the following way. First, we choose a sequence $k_{n}$ of positive integers and a number $\theta \in(0, \infty)$ such that $k_{n}=\theta \sqrt{n}+\mathrm{o}\left(n^{1 / 4}\right)$ as $n \rightarrow \infty$. We also choose a continuous function $g:[0,1] \rightarrow \mathbb{R}$ which is piecewise $C^{1}$ with a piecewise Lipschitz derivative $g^{\prime}$ and satisfies $g(0)=g(1)=0$ and $\int_{0}^{1} g(x)^{2} \mathrm{~d} x>0$. After that, for any $d$-dimensional stochastic process $V=\left(V^{1}, \ldots, V^{d}\right)$ we define the quantity

$$
\begin{equation*}
\bar{V}_{i}^{k}=\sum_{p=1}^{k_{n}-1} g\left(\frac{p}{k_{n}}\right)\left(V_{\tau_{i+p}^{k}}^{k}-V_{\tau_{i+p-1}^{k}}^{k}\right) \tag{2.3}
\end{equation*}
$$

and set $\bar{V}_{i}=\left(\bar{V}_{i}^{1}, \ldots, \bar{V}_{i}^{d}\right)^{*}$ (hereafter an asterisk denotes the transpose of a matrix). The MRC estimator is defined by

$$
\operatorname{MRC}[Y]_{t}^{n}=\frac{1}{\psi_{2} k_{n}} \sum_{i=0}^{N_{t}^{n}-k_{n}+1} \bar{Y}_{i}\left(\bar{Y}_{i}\right)^{*}-\frac{\psi_{1}}{2 \psi_{2} k_{n}^{2}}[Y, Y]_{t}^{n},
$$

where $N_{t}^{n}=\max \left\{p: T_{p} \leq t\right\}, \psi_{1}=\int_{0}^{1} g^{\prime}(x)^{2} \mathrm{~d} x, \psi_{2}=\int_{0}^{1} g(x)^{2} \mathrm{~d} x$ and

$$
[Y, Y]_{t}^{n}=\sum_{p=1}^{N_{t}^{n}} \Delta_{p} Y\left(\Delta_{p} Y\right)^{*}, \quad \Delta_{p} Y=\left(Y_{\tau_{p}^{1}}^{1}-Y_{\tau_{p-1}^{1}}^{1}, \ldots, Y_{\tau_{p}^{d}}^{d}-Y_{\tau_{p-1}^{d}}^{d}\right)^{*}
$$

for each $t \in \mathbb{R}_{+}$. Here, we set $\sum_{i=p}^{q} \equiv 0$ if $p>q$ by convention. In the synchronous and equidistant sampling setting, the asymptotic distribution of the MRC estimator has been derived in [29]
(see Theorem 4.6 of that paper, and see also Section 4 of [21]). Our purpose is to develop the asymptotic distribution of the MRC estimator in the situation where the observation times are possibly irregular, non-synchronous and endogenous.

## 3. Main result

### 3.1. Notation

In this subsection, some notation is introduced in order to state our main result. First we introduce notation appearing in the assumptions stated in the next subsection. We write $X^{n} \xrightarrow{\text { ucp }} X$ for processes $X^{n}$ and $X$ to express shortly that $\sup _{0 \leq t \leq T}\left|X_{t}^{n}-X_{t}\right| \rightarrow^{p} 0$ for any $T>0$. $\varpi$ denotes some (fixed) positive constant. We denote by $\left(\mathcal{G}_{t}^{(0)}\right)$ (resp., $\left(\mathcal{G}_{t}\right)$ ) the smallest filtration containing $\left(\mathcal{F}_{t}^{(0)}\right)$ (resp., $\left.\left(\mathcal{F}_{t}\right)\right)$ such that $\mathcal{G}_{0}^{(0)}$ (resp., $\mathcal{G}_{0}$ ) contains the $\sigma$-field generated by $\mu$, that is, the $\sigma$-field generated by all the variables $\mu(A)$, where $A$ ranges all measurable subsets of $\mathbb{R}_{+} \times E$.

Next, we introduce some quantities appearing in the representation of the asymptotic variance of the estimator. We set $\Sigma_{s}=\sigma_{s} \sigma_{s}^{*}$ for each $s \in \mathbb{R}_{+}$, that is, $\Sigma$ denotes the diffusion coefficient matrix process. We denote by $\Upsilon_{t}$ the covariance matrix of $\epsilon_{t}$, that is, $\Upsilon_{t}(\cdot)=\int u u^{*} Q_{t}(\cdot, \mathrm{~d} u)$ (we will assume the existence of the second moment of the noise later, so this matrix always exists). For any real-valued bounded measurable functions $u, v$ on $[0,1]$, we define the function $\phi_{u, v}$ on $[0,1]$ by $\phi_{u, v}(y)=\int_{y}^{1} u(x-y) v(x) \mathrm{d} x$. Then, we put

$$
\Phi_{22}=\int_{0}^{1} \phi_{g, g}(y)^{2} \mathrm{~d} y, \quad \Phi_{12}=\int_{0}^{1} \phi_{g, g}(y) \phi_{g^{\prime}, g^{\prime}}(y) \mathrm{d} y, \quad \Phi_{11}=\int_{0}^{1} \phi_{g^{\prime}, g^{\prime}}(y)^{2} \mathrm{~d} y .
$$

On the other hand, for any $k, l=1, \ldots, d$ we define the process $\mathfrak{J}^{k l}$ by

$$
\mathfrak{J}_{s}^{k l}=\Delta X_{s}^{k} \Delta X_{s}^{l}\left\{\Phi_{22} \theta\left(\Sigma_{s-}^{k l} G_{s-}+\Sigma_{s}^{k l} G_{s}\right)+\frac{\Phi_{12}}{\theta}\left(\Upsilon_{s-}^{k l} \chi_{s-}^{k l}+\Upsilon_{s}^{k l} \chi_{s}^{k l}\right)\right\} .
$$

Remark 3.1 (Properties of $\phi_{u, v}$ ). We will use the following properties of $\phi_{u, v}$ : first, for any realvalued bounded measurable function $u$ on $[0,1], \phi_{u, u}$ is non-negative. In fact, setting $u(x)=0$ for $x \notin[0,1]$, we have $\phi_{u, u}(y)=\int_{-\infty}^{\infty} u(-(y-x)) u(x) \mathrm{d} x$ for all $y \in[0,1]$, and we can extend the domain of $\phi_{u, u}$ to the whole real line using this expression. Then, denoting by $\hat{f}$ the Fourier transform of a function $f$ on $\mathbb{R}$, we have $\hat{\phi}_{u, u}=|\hat{u}|^{2} \geq 0$. Hence, $\phi_{u, u}$ is a positive definite function and, in particular, $\phi_{u, u}(y) \geq 0$ for all $y \in \mathbb{R}$. Next, we can easily check that $\phi_{g, g}^{\prime}=$ $\phi_{g, g^{\prime}}=-\phi_{g^{\prime}, g}$ and $\phi_{g, g}^{\prime \prime}=-\phi_{g^{\prime}, g^{\prime}}$. In particular, $\Phi_{12}=\int_{0}^{1} \phi_{g^{\prime}, g}(y)^{2} \mathrm{~d} y$ due to integration by parts.

### 3.2. Assumptions

We impose the following condition on the sampling schemes $\left(T_{p}\right)_{p \geq 0}$ and $\left(\tau_{p}^{k}\right)_{p \geq 0}(k=$ $1, \ldots, d)$ :
[A1] $\left(T_{p}\right)_{p \geq 0}$ and $\left(\tau_{p}^{k}\right)_{p \geq 0}(k=1, \ldots, d)$ are sequences of $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping times and satisfy (2.2). It also holds that

$$
\begin{equation*}
r_{n}(t):=\sup _{p \geq 0}\left(T_{p} \wedge t-T_{p-1} \wedge t\right)=\mathrm{o}_{p}\left(n^{-\xi}\right) \tag{3.1}
\end{equation*}
$$

as $n \rightarrow \infty$ (note that $T_{-1}=0$ by convention) for every $t>0$ and every $\xi \in(0,1)$. Moreover, for each $n$ we have a $\left(\mathcal{G}_{t}^{(0)}\right)$-progressively measurable positive-valued process $G_{t}^{n}$, a $\left(\mathcal{G}_{t}^{(0)}\right)$ progressively measurable $[0,1]^{d} \otimes[0,1]^{d}$-valued process $\chi_{t}^{n}=\left(\chi_{t}^{n, k l}\right)_{1 \leq k, l \leq d}$ and a random subset $\mathcal{N}^{n}$ of $\mathbb{Z}_{+}$satisfying the following conditions:
(i) $\left\{(\omega, p) \in \Omega \times \mathbb{Z}_{+}: p \in \mathcal{N}^{n}(\omega)\right\}$ is a measurable set of $\Omega \times \mathbb{Z}_{+}$. Moreover, there is a constant $\kappa \in\left(0, \frac{1}{2}\right)$ such that $\#\left(\mathcal{N}^{n} \cap\left\{p: T_{p} \leq t\right\}\right)=\mathrm{O}_{p}\left(n^{\kappa}\right)$ as $n \rightarrow \infty$ for every $t>0$.
(ii) $E\left[n\left(T_{p+1}-T_{p}\right) \mid \mathcal{G}_{T_{p}}^{(0)}\right]=G_{T_{p}}^{n}$ and $E\left[1_{\left\{\tau_{p+1}^{k}=\tau_{p+1}^{l}\right\}} \mid \mathcal{G}_{T_{p}}^{(0)}\right]=\chi_{T_{p}}^{n, k l}$ for every $n$, every $p \in$ $\mathbb{Z}_{+}-\mathcal{N}^{n}$ and any $k, l=1, \ldots, d$.
(iii) There is a cadlag $\left(\mathcal{F}_{t}^{(0)}\right)$-adapted positive valued process $G$ such that:
(iii-a) $n^{\sigma}\left(G^{n}-G\right) \xrightarrow{\text { ucp }} 0$,
(iii-b) $G_{t-}>0$ for every $t>0$,
(iii-c) $G$ is an Itô semimartingale of the form

$$
G_{t}=G_{0}+\int_{0}^{t} \widehat{b}_{s} \mathrm{~d} s+\int_{0}^{t} \widehat{\sigma}_{s} \mathrm{~d} W_{s}+\left(\widehat{\delta} 1_{\{|\widehat{\delta}| \leq 1\}}\right) \star(\mu-v)_{t}+\left(\widehat{\delta} 1_{\{|\widehat{\delta}|>1\}}\right) \star \mu_{t}
$$

where $\widehat{b}_{s}$ is a locally bounded and $\left(\mathcal{F}_{t}^{(0)}\right)$-progressively measurable real-valued process, $\widehat{\sigma}_{s}$ is a cadlag $\left(\mathcal{F}_{t}^{(0)}\right)$-adapted $\mathbb{R} \otimes \mathbb{R}^{d^{\prime}}$-valued process, and $\widehat{\delta}$ is an $\left(\mathcal{F}_{t}^{(0)}\right)$-predictable realvalued function on $\Omega^{(0)} \times \mathbb{R}_{+} \times E$ such that there is a sequence $\left(\widehat{\rho}_{j}\right)$ of $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping times increasing to infinity and, for each $j$, a deterministic non-negative function $\widehat{\gamma}_{j}$ on $E$ satisfying $\int \widehat{\gamma}_{j}(z)^{2} \wedge 1 \lambda(\mathrm{~d} z)<\infty$ and $\left|\widehat{\delta}\left(\omega^{(0)}, t, z\right)\right| \leq \widehat{\gamma}_{j}(z)$ for all $\left(\omega^{(0)}, t, z\right)$ with $t \leq \widehat{\rho}_{j}\left(\omega^{(0)}\right)$.
(iv) There is a cadlag $\left(\mathcal{F}_{t}^{(0)}\right)$-adapted $[0,1]^{d} \otimes[0,1]^{d}$-valued process $\chi$ such that $n^{\varpi}\left(\chi^{n}-\right.$ $\chi) \xrightarrow{\text { ucp }} 0$ as $n \rightarrow \infty$. Furthermore, for each $j \in \mathbb{N}$ we have a cadlag $\left(\mathcal{F}_{t}^{(0)}\right)$-adapted $[0,1]^{d} \otimes[0,1]^{d}$-valued process $\chi(j)$, an $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping time $\check{\rho}_{j}$, and a constant $\check{\Lambda}_{j}$ such that $\check{\rho}_{j} \uparrow \infty$ as $j \rightarrow \infty$ and $\chi\left(\omega^{(0)}\right)_{t}=\chi(j)\left(\omega^{(0)}\right)_{t}$ if $t<\check{\rho}_{j}\left(\omega^{(0)}\right)$ and

$$
E\left[\left\|\chi(j)_{t_{1}}-\chi(j)_{t_{2}}\right\|^{2} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right] \leq \check{\Lambda}_{j} E\left[\left|t_{1}-t_{2}\right|^{\sigma} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right]
$$

for every $j$ and any $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping times $t_{1}$ and $t_{2}$ bounded by $j$.
Remark 3.2. (i) The assumptions on ( $T_{p}$ ) are motivated by the concept of the restricted discretization scheme discussed in detail in Chapter 14 of [31]. In fact, suppose that $T_{p}$ 's are of the form

$$
T_{p}=T_{p-1}+\theta_{T_{p-1}}^{n} \varepsilon(n, p), \quad p=1,2, \ldots
$$

where $\theta^{n}$ is a cadlag $\left(\mathcal{F}_{t}^{(0)}\right)$-adapted process, $(\varepsilon(n, p))_{p \geq 1}$ is a sequence of i.i.d. positive variables independent of $b, \sigma, \delta, W, \mu$, and such that $E[\varepsilon(n, p)]=1$ and $E\left[\varepsilon(n, p)^{r}\right]<\infty$ for every $r>0$, and $T_{0}=0$. By constructing the filtration $\left(\mathcal{F}_{t}^{(0)}\right)$ suitably, we may assume that $\epsilon(n, p)$ is independent of $\mathcal{F}_{T_{p-1}}^{(0)}$ for all $n, p$. Then we have [A1](i)-(ii) regarding $G^{n}$ while we set $\mathcal{N}^{n}=\varnothing$ and $G^{n}=n \theta^{n}$. In this case [A1](iii) corresponds to (a weaker version of) assumption (E) from [31], and (3.1) follows from Lemma 14.1.5 of [31]. Unlike their setting, however, our assumption does not rule out the dependence between $\varepsilon(n, p)$ 's and $X$ (see, e.g., Example 4.1 in the next section). The importance of such dependence has recently been emphasized in econometric literature; see, for example, [49].
(ii) The assumptions on the quantities $1_{\left\{\tau_{p}^{k}=\tau_{p}^{l}\right\}}$ are necessary for the treatment of the $\left(\mathcal{F}^{(0)}\right.$ conditional) covariance between $\epsilon_{\tau_{p}^{k}}^{k}$ and $\epsilon_{\tau_{p}^{l}}^{l}$, which is given by $\Upsilon_{\tau_{p}^{k}}^{k l} 1_{\left\{\tau_{p}^{k}=\tau_{p}^{l}\right\}}$ (a similar kind of assumption also appears in [10] due to the same reason as ours). Therefore, those assumptions can be dropped when $\Upsilon^{k l} \equiv 0$ if $k \neq l$; this is often assumed in the literature on the covariance estimation of non-synchronously observed semimartingales with noise. The quantity $\chi^{n}$ measures the degree of the non-synchronicity, and $\chi_{s}^{n}$ is a matrix all of whose components are equal to 1 in the synchronous case while it is an identity matrix in the completely non-synchronous case. Hence, [A1](iv) is satisfied in these two extreme cases.
(iii) The possibility of the set $\mathcal{N}^{n}$ being non-empty excludes the following trivial exception of [A1] with $\mathcal{N}^{n}$ being empty: if $T_{0}=\log n / n$ and $T_{p}=T_{p-1}+1 / n$ for $p \geq 1$, [A1] with $\mathcal{N}^{n}=\varnothing$ is not satisfied because $G_{T_{0}}^{n} \rightarrow \infty$ as $n \rightarrow \infty$. This assumption is also useful to ensure the stability under the localization used in the proof; see Lemma 6.1.
(iv) The fact that we consider the conditional expected durations given $\mathcal{G}_{T_{p}}^{(0)}$,s instead of $\mathcal{F}_{T_{p}}^{(0)}$ rules out some dependence between the sampling schemes and the jumps of the observed process. For example, if $\mu$ is a jump measure of a one-dimensional Lévy process (i.e., $E=\mathbb{R}$ ) and $T_{p}$ 's are of the form $T_{p}=\inf \left\{t>T_{p-1}:\left|\int_{T_{p-1}}^{t} \int_{|z| \leq 1} z(\mu-v)(\mathrm{d} s, \mathrm{~d} z)\right|>\eta_{n}\right\}$ for $p=1,2, \ldots$ and for some appropriate sequence $\left(\eta_{n}\right)_{n \geq 1}$ of positive numbers, then [A1] obviously fails because $T_{p}$ 's are $\mathcal{G}_{0}^{(0)}$-measurable (this type of sampling scheme is well studied in [50]). On the other hand, it still allows the presence of the instantaneous causality between the sampling schemes and the jumps: see Example 4.2.
(v) Under [A1], it holds that

$$
\begin{equation*}
\frac{1}{n} N_{t}^{n} \rightarrow^{p} \int_{0}^{t} \frac{1}{G_{s}} \mathrm{~d} s \tag{3.2}
\end{equation*}
$$

as $n \rightarrow \infty$ for every $t \in \mathbb{R}_{+}$(see Section 6.1 of [36] for the proof). In particular, [A1] ensures that the parameter $n$ controls the magnitude of the number of observations.

We impose the following structural assumption on the latent process $X$ :
[A2] The volatility process $\sigma$ is an Itô semimartingale of the form

$$
\sigma_{t}=\sigma_{0}+\int_{0}^{t} \widetilde{b}_{s} \mathrm{~d} s+\int_{0}^{t} \widetilde{\sigma}_{s} \mathrm{~d} W_{s}+\left(\widetilde{\delta} 1_{\{|\tilde{\delta}| \leq 1\}}\right) \star(\mu-v)_{t}+\left(\widetilde{\delta} 1_{\{|\tilde{\delta}|>1\}}\right) \star \mu_{t}
$$

where $\widetilde{b}_{s}$ is a locally bounded and $\left(\mathcal{F}_{t}^{(0)}\right)$-progressively measurable $\mathbb{R}^{d} \otimes \mathbb{R}^{d^{\prime}}$-valued process, $\widetilde{\sigma}_{s}$ is a cadlag $\left(\mathcal{F}_{t}^{(0)}\right)$-adapted $\mathbb{R}^{d} \otimes \mathbb{R}^{d^{\prime}} \otimes \mathbb{R}^{d^{\prime}}$-valued process, and $\widetilde{\delta}$ is an $\left(\mathcal{F}_{t}^{(0)}\right)$-predictable $\mathbb{R}^{d} \otimes \mathbb{R}^{d^{\prime}}$-valued function on $\Omega^{(0)} \times \mathbb{R}_{+} \times E$.

Moreover, for each $j$ there is an $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping time $\rho_{j}$, a bounded $\left(\mathcal{F}_{t}^{(0)}\right)$-progressively measurable $\mathbb{R}^{d}$-valued process $b(j)_{s}$, a deterministic non-negative function $\gamma_{j}$ on $E$, and a constant $\Lambda_{j}$ such that $\rho_{j} \uparrow \infty$ as $j \rightarrow \infty$ and, for each $j$,
(i) $b\left(\omega^{(0)}\right)_{s}=b(j)\left(\omega^{(0)}\right)_{s}$ if $s<\rho_{j}\left(\omega^{(0)}\right)$,
(ii) $E\left[\left\|b(j)_{t_{1}}-b(j)_{t_{2}}\right\|^{2} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right] \leq \Lambda_{j} E\left[\left|t_{1}-t_{2}\right|^{\sigma} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right]$ for any $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping times $t_{1}$ and $t_{2}$ bounded by $j$,
(iii) $\int\left\{\gamma_{j}(z)^{2} \wedge 1\right\} \lambda(\mathrm{d} z)<\infty$ and $\left\|\delta\left(\omega^{(0)}, t, z\right)\right\| \vee\left\|\widetilde{\delta}\left(\omega^{(0)}, t, z\right)\right\| \leq \gamma_{j}(z)$ for all $\left(\omega^{(0)}, t, z\right)$ with $t \leq \rho_{j}\left(\omega^{(0)}\right)$,
(iv) $E\left[\left\|\delta\left(t_{1} \wedge \rho_{j}, z\right)-\delta\left(t_{2} \wedge \rho_{j}, z\right)\right\|^{2} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right] \leq \Lambda_{j} \gamma_{j}(z)^{2} E\left[\left|t_{1}-t_{2}\right|^{\sigma} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right]$ for any $\left(\mathcal{F}_{t}^{(0)}\right)$ stopping times $t_{1}$ and $t_{2}$ bounded by $j$.

Remark 3.3. An [A2] type assumption is commonly used in the literature of power variations (see, e.g., [31]), except for assumptions (ii) and (iv), that is, continuity assumptions on the drift and the jump coefficient. Such assumptions are necessary for the treatment of the irregularity and the non-synchronicity of the observation times as in [24].

We also impose the following regularity condition on the noise process:
[A3] There is a constant $\Gamma>4$ and a sequence $\left(\rho_{j}^{\prime}\right)_{j \geq 1}$ of $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping times increasing to infinity such that

$$
\sup _{\omega^{(0)} \in \Omega^{(0)}, t<\rho_{j}^{\prime}\left(\omega^{(0)}\right)} \int\|z\|^{\Gamma} Q_{t}\left(\omega^{(0)}, \mathrm{d} z\right)<\infty .
$$

Moreover, for each $j$ there is a bounded cadlag $\left(\mathcal{F}_{t}^{(0)}\right)$-adapted $\mathbb{R}^{d} \otimes \mathbb{R}^{d}$-valued process $\Upsilon(j)_{t}$ and a constant $\Lambda_{j}^{\prime}$ such that:
(i) $\Upsilon(j)\left(\omega^{(0)}\right)_{t}=\Upsilon\left(\omega^{(0)}\right)_{t}$ if $t<\rho_{j}^{\prime}\left(\omega^{(0)}\right)$,
(ii) $E\left[\left\|\Upsilon(j)_{t_{1}}-\Upsilon(j)_{t_{2}}\right\|^{2} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right] \leq \Lambda_{j}^{\prime} E\left[\left|t_{1}-t_{2}\right|^{\Phi} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right]$ for any $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping times $t_{1}$ and $t_{2}$ bounded by $j$.

Remark 3.4. The locally boundedness of the moment process of the noise is used for verifying a Lyapunov type condition for central limit theorems and proving the negligibility of the edge effect. The continuity assumption of the covariance matrix process of the noise is necessary due to the same reason as for [A2]. If the noise is assumed to be i.i.d. and independent of $\mathcal{F}^{(0)}$, [A3] simply means the $\Gamma$ th moment of the noise is finite for some $\Gamma>4$.

Finally, we introduce the following technical condition to avoid some measure-theoretic problems:
[A4] (i) A regular conditional probability of $P^{(0)}$ given $\mathcal{H}$ exists for any sub- $\sigma$-field $\mathcal{H}$ of $\mathcal{F}^{(0)}$.
(ii) The process $\left(Q_{t}(\cdot, A)\right)_{t \geq 0}$ is $\left(\mathcal{F}_{t}^{(0)}\right)$-progressively measurable for any Borel set $A$ of $\mathbb{R}^{d}$.

Remark 3.5. (i) [A4](i) is satisfied, for example, when $\left(\Omega^{(0)}, \mathcal{F}^{(0)}\right)$ is a standard measurable space, that is, it is Borel isomorphic to some Polish space (see, e.g., Theorem I-3.1 of [26]). In fact, this assumption is not restrictive for applications.
(ii) [A4](ii) is satisfied, for example, when $Q_{t} \equiv Q$ for some probability measure $Q$ on $\mathbb{R}^{d}$, that is, the noise is modeled by an i.i.d. sequence. Another example is the case where $Q_{t}\left(\omega^{(0)}, \cdot\right)$ has a density of the form $f\left(\cdot, X_{t}\left(\omega^{(0)}\right)\right.$ ), where $f$ is a measurable function on $\mathbb{R}^{d} \times \mathbb{R}^{d}$ into [0, 1] such that $\int_{\mathbb{R}^{d}} f(x, \theta) \mathrm{d} x=1$ for every $\theta \in \mathbb{R}^{d}$. Example 16.1.5 of [31] is encompassed with this type of model. Thus, this assumption also seems to be unrestrictive for applications.

### 3.3. Result

To state the main result, we need the notion of stable convergence as common in this area. For each $n \geq 1$, let $X^{n}$ be a random variable which is defined on $\mathcal{B}$ and takes values in a Polish space $S$. The variables $X^{n}$ are said to converge stably in law to an $S$-valued random variable $X$ defined on an extension of $\mathcal{B}^{(0)}$ if $E\left[U f\left(X^{n}\right)\right] \rightarrow \widetilde{E}[U f(X)]$ for any $\mathcal{F}^{(0)}$-measurable bounded random variable $U$ and any bounded continuous function $f$ on $S$, where $\widetilde{E}$ denotes the expectation with respect to the probability measure of the extension. We then write $X^{n} \rightarrow^{d_{s}} X$. Note that we need a slightly generalized definition of stable convergence described at the end of Section 2.2.1 of [31] because $\mathcal{B}$ changes as $n$ varies. The most important property of stable convergence is the following: if the real-valued variables $V_{n}$ defined on $\mathcal{B}$ converge in probability to a variable $V$ defined on $\mathcal{B}^{(0)}$, then $X^{n} \rightarrow^{d_{s}} X$ implies that $\left(X^{n}, V_{n}\right) \rightarrow^{d_{s}}(X, V)$ for the product topology on the space $S \times \mathbb{R}$.

Theorem 3.1. Suppose that [A1]-[A4] are satisfied. Then

$$
n^{1 / 4}\left(\operatorname{MRC}[Y]_{t}^{n}-[X, X]_{t}\right) \rightarrow^{d_{s}} \mathcal{W}_{t}+\mathcal{Z}_{t}
$$

as $n \rightarrow \infty$ for any $t>0$, where $\mathcal{W}$ and $\mathcal{Z}$ are $\mathbb{R}^{d} \otimes \mathbb{R}^{d}$-valued processes defined on an extension of $\mathcal{B}^{(0)}$, which conditionally on $\mathcal{F}^{(0)}$ are mutually independent, centered Gaussian with independent increments, the first one being continuous and the second one being purely discontinuous, and with (conditional) covariances

$$
\begin{align*}
\widetilde{E}[ & \left.\mathcal{W}_{t}^{k l} \mathcal{W}_{t}^{k^{\prime^{\prime}}} \mid \mathcal{F}^{(0)}\right] \\
=\frac{2}{\psi_{2}^{2}} \int_{0}^{t}[ & {\left[\Phi_{22} \theta\left\{\Sigma_{s}^{k k^{\prime}} \Sigma_{s}^{l l^{\prime}}+\Sigma_{s}^{k l^{\prime}} \Sigma_{s}^{l k^{\prime}}\right\} G_{s}\right.}  \tag{3.3}\\
& +\frac{\Phi_{11}}{\theta^{3}}\left\{\Upsilon_{s}^{k k^{\prime}} \chi_{s}^{k k^{\prime}} \Upsilon_{s}^{l l^{\prime}} \chi_{s}^{l l^{\prime}}+\Upsilon_{s}^{k l^{\prime}} \chi_{s}^{k l^{\prime}} \Upsilon_{s}^{l k^{\prime}} \chi_{s}^{l k^{\prime}}\right\} \frac{1}{G_{s}} \\
& \left.+\frac{\Phi_{12}}{\theta}\left\{\Sigma_{s}^{k k^{\prime}} \Upsilon_{s}^{l l^{\prime}} \chi_{s}^{l l^{\prime}}+\Sigma_{s}^{l k^{\prime}} \Upsilon_{s}^{k l^{\prime}} \chi_{s}^{k l^{\prime}}+\Sigma_{s}^{l l^{\prime}} \Upsilon_{s}^{k k^{\prime}} \chi_{s}^{k k^{\prime}}+\Sigma_{s}^{k l^{\prime}} \Upsilon_{s}^{l k^{\prime}} \chi_{s}^{l k^{\prime}}\right\}\right] \mathrm{d} s
\end{align*}
$$

and

$$
\begin{equation*}
\widetilde{E}\left[\mathcal{Z}_{t}^{k l} \mathcal{Z}_{t}^{k^{\prime} l^{\prime}} \mid \mathcal{F}^{(0)}\right]=\frac{1}{\psi_{2}^{2}} \sum_{s \leq t}\left(\mathfrak{J}_{s}^{k k^{\prime}}+\mathfrak{J}_{s}^{k l^{\prime}}+\mathfrak{J}_{s}^{l k^{\prime}}+\mathfrak{J}_{s}^{l l^{\prime}}\right) \tag{3.4}
\end{equation*}
$$

Here, $\widetilde{E}$ denotes the expectation with respect to the probability measure of the extension.
When further $X$ is continuous, the processes $n^{1 / 4}\left(\operatorname{MRC}[Y]^{n}-[X, X]\right)$ converge stably in law to the process $\mathcal{W}$ for the Skorokhod topology.

Remark 3.6. The above theorem shows that the observation times' effect on the asymptotic distribution of the MRC estimator is only through the asymptotic conditional expected duration process $G$ and the limiting process $\chi$ measuring the degree of the non-synchronicity. As was indicated in Remark 3.2(ii), $\chi$ simply reflects the covariance structure of the noise process, while $G$ naturally affects the asymptotic distribution of the estimator because it links with the (spot) sampling frequency, as seen from (3.2). Consequently, the irregularity and the endogeneity of the observation times have no impact on the asymptotic distribution of the estimator.

Remark 3.7. In the proof of the theorem, it plays a key role to replace the duration ( $T_{p+1}-T_{p}$ ) with its conditional expectation $G_{T_{p}}^{n}$. Such replacement is possible because our estimator contains a local averaging procedure (2.3). More formally, this procedure makes it possible to apply a standard martingale argument described in Lemma 2.3 of [18] to the durations. The benefits of this fact appear in the treatments of the irregularity and the endogeneity of the observation times in Lemmas 6.4 and 6.7. Also, this is why the higher (conditional) moments of the durations do not affect the asymptotic distribution of the estimator.

Remark 3.8 (Covariance structure of $\mathcal{W}_{t}$ ). It is convenient to observe that the covariance structure of $\mathcal{W}_{t}$ is analogous to the asymptotic covariance of the realized covariance in a standard setting. For this purpose, in the following we use some concepts from matrix algebra found in, for example, [25]. For each $s \in \mathbb{R}_{+}$, we denote by $\widetilde{\Upsilon}_{s}$ the Hadamard product of $\Upsilon_{s}$ and $\chi_{s}$, that is, $\widetilde{\Upsilon}_{s}^{k l}=\Upsilon_{s}^{k l} \chi_{s}^{k l}$ for $k, l=1, \ldots, d$, and set $\bar{\Sigma}_{s}(y)=\frac{\sqrt{2}}{\psi_{2}}\left\{\phi_{g, g}(y) \theta^{1 / 2} \Sigma_{s} \sqrt{G_{s}}+\right.$ $\left.\phi_{g^{\prime}, g^{\prime}}(y) \theta^{-3 / 2} \widetilde{\Upsilon}_{s} / \sqrt{G_{s}}\right\}$. Since both $\Upsilon_{s}$ and $\chi_{s}$ is positive semidefinite, so is $\widetilde{\Upsilon}_{s}$ due to the Schur product theorem, that is, Theorem 5.2 .1 of [25] (note that the positive semidefiniteness of $\chi_{s}$ can be checked directly using the fact that $\chi_{s}^{k k}=1$ and $0 \leq \chi_{s}^{k l} \leq 1$ for any $\left.k, l\right)$. Therefore, $\bar{\Sigma}_{s}(y)$ is positive semidefinite as well because both $\phi_{g, g}$ and $\phi_{g^{\prime}, g^{\prime}}$ are non-negative (see Remark 3.1). Then the left-hand side of (3.3) can be rewritten as

$$
\int_{0}^{t}\left[\int_{0}^{1}\left\{\bar{\Sigma}_{s}(y)^{k k^{\prime}} \bar{\Sigma}_{s}(y)^{l l^{\prime}}+\bar{\Sigma}_{s}(y)^{k l^{\prime}} \bar{\Sigma}_{s}(y)^{l k^{\prime}}\right\} \mathrm{d} y\right] \mathrm{d} s
$$

The integrand of the above expression is nothing but the $\mathcal{F}^{(0)}$-conditional covariance between the $(k, l)$ th and $\left(k^{\prime}, l^{\prime}\right)$ th entries of the variable $\bar{\Sigma}_{s}(y)^{1 / 2} \zeta\left(\bar{\Sigma}_{s}(y)^{1 / 2} \zeta\right)^{*}$, where $\zeta$ is a $d$-dimensional standard normal variable independent of $\mathcal{F}^{(0)}$. In other words, vec $\left(\mathcal{W}_{t}\right)$ is centered Gaussian with covariance matrix $\int_{0}^{1} \mathfrak{S}_{s} \mathrm{~d} s$, where $\mathfrak{S}_{s}=\int_{0}^{1}\left(\bar{\Sigma}_{s}(y) \otimes \bar{\Sigma}_{s}(y)\right) \operatorname{Cov}\left[\operatorname{vec}\left(\zeta \zeta^{*}\right)\right] \mathrm{d} y$ and, vec and $\otimes$ denote the vec-operator and the Kronecker product of matrices, respectively (cf. Section 2.2
of [9]). In particular, the process $\mathfrak{S}_{s}$ is cadlag, $\left(\mathcal{F}_{t}^{(0)}\right)$-adapted and takes values in the set of $d \times d$ positive semidefinite matrices, hence we can construct the process $\mathcal{W}$ stated as in the theorem by Proposition 4.1.2 of [31]. More precisely, $\mathcal{W}$ can be realized as $\operatorname{vec}\left(\mathcal{W}_{t}\right)=\int_{0}^{t} \mathfrak{S}_{s}^{1 / 2} \mathrm{~d} W_{s}^{\prime}$, where $W^{\prime}$ is a $d^{2}$-dimensional standard Brownian motion defined on an extension of $\mathcal{B}^{(0)}$ and independent of $\mathcal{F}^{(0)}$.

Note that the Fisher information matrix for covariance matrix estimation of a multivariate diffusion process from non-synchronous and noisy observations is not analogous to that for a pure diffusion setting; see Section 2.2 of [9] for details.

Remark 3.9 (Covariance structure of $\mathcal{Z}_{t}$ ). $\mathcal{Z}_{t}$ apparently has an analogous covariance structure to the asymptotic covariance of the realized covariance due to jumps in the regular sampling case (cf. equation (5.4.4) of [31]), and it can be realized as follows. Set $A_{m}=\{z: \gamma(z)>1 / m\}$ for each $m \in \mathbb{N}$, and denote by $(S(m, j))_{j \geq 1}$ the successive jump times of the Poisson process $1_{A_{m} \backslash A_{m-1}} \star \mu$. Let $\left(S_{r}\right)_{r \geq 1}$ be a reordering of the double sequence $(S(m, j))$. Suppose that sequences $\left(\Psi_{r-}\right)_{r \geq 1}$ and $\left(\Psi_{r+}\right)_{r \geq 1}$ of i.i.d. standard $d^{\prime}$-dimensional normal variables and sequences $\left(\Psi_{r-}^{\prime}\right)_{r \geq 1}$ and $\left(\Psi_{r+}^{\prime}\right)_{r \geq 1}$ of i.i.d. standard $d$-dimensional normal variables are defined on an extension of $\mathcal{B}^{(0)}$ and that all of them are mutually independent and independent of $\mathcal{F}^{(0)}$. Now, the variable $\widetilde{\Upsilon}_{s}$ defined in Remark 3.8 is positive semidefinite, it admits the (positive semidefinite) square root $\widetilde{v}_{s}:=\Upsilon_{s}^{1 / 2}$. Since the process $\widetilde{\Upsilon}$ is cadlag and $\left(\mathcal{F}_{t}^{(0)}\right)$-adapted, so is $\widetilde{v}_{s}$. Then $\mathcal{Z}$ is realized as $\mathcal{Z}_{t}=\sum_{r: S_{r} \leq t}\left(\mathfrak{Z}_{r}+\mathfrak{Z}_{r}^{*}\right)$, where

$$
\mathfrak{Z}_{r}=\frac{1}{\psi_{2}} \Delta X_{S_{r}}\left\{\sqrt{\Phi_{22} \theta}\left(\sigma_{S_{r}-} \sqrt{G_{S_{r}-}} \Psi_{r-}+\sigma_{S_{r}} \sqrt{G_{S_{r}}} \Psi_{r+}\right)+\sqrt{\frac{\Phi_{12}}{\theta}}\left(\widetilde{v}_{S_{r}-} \Psi_{r-}^{\prime}+\widetilde{v}_{S_{r}} \Psi_{r+}^{\prime}\right)\right\}^{*}
$$

This is indeed the desired one; see Proposition 4.1.4 of [31].
Remark 3.10 (Comparison with a pure semimartingale setting). It would be interesting to observe how our result is different from [12]'s one in a pure semimartingale setting. For simplicity, we focus on the univariate case, that is, we assume that $d=d^{\prime}=1$, and assume that $T_{p}=t_{p}^{1}$ for every $p$ for notational simplicity. Now let us recall their result briefly. Suppose that $b, \sigma$ and $\delta$ are continuous. Suppose also that the sequence $\left(T_{p}\right)$ is independent of $b, \sigma, \delta, W, \mu$ and satisfies (3.1). Then, according to Theorem 2 of [12], for any $t>0$ we have the following convergence:

$$
\begin{equation*}
\sqrt{n}\left([X, X]_{t}^{n}-[X, X]_{t}\right) \rightarrow^{d_{s}} \sqrt{2} \int_{0}^{t} \sigma_{s}^{2} \sqrt{H^{\prime}(s)} \mathrm{d} W_{s}^{\prime}+2 \sum_{r: S_{r} \leq t} \Delta X_{S_{r}} \sigma_{S_{r}} \sqrt{\eta\left(S_{r}\right)} \Psi_{r}, \tag{3.5}
\end{equation*}
$$

where $W^{\prime}$ is a standard Brownian motion, $H$ is a (possibly random) $C^{1}$ function such that $n \sum_{p: T_{p} \leq t}\left(T_{p}-T_{p-1}\right)^{2} \rightarrow^{p} H(t)$ for every $t \in \mathbb{R}_{+}$(the existence is assumed), $\left(S_{r}\right)_{r \geq 1}$ is a sequence of stopping times exhausting the jumps of $X,\left(\Psi_{r}\right)_{r \geq 1}$ is a sequence of i.i.d. standard normal variables, and $(\eta(t))_{t \in \mathbb{R}_{+}}$is a family of independent random variables with uniformly bounded first moments, and such that the processes $\left(n\left(T_{+}(t)-T_{-}(t)\right)\right)_{t \in \mathbb{R}_{+}}$converge finite-dimensionally in law to $(\eta(t))_{t \in \mathbb{R}_{+}}$(the existence is assumed, and this condition can
be weakened; see Assumption 2 of [12] for details). Here, $T_{+}(t)=\min \left\{T_{p}: T_{p} \geq t\right\}$ and $T_{-}(t)=\max \left\{T_{p}: T_{p} \leq t\right\}$ for any $t \in \mathbb{R}_{+}$and $W^{\prime},\left(\Psi_{r}\right)$ and $(\eta(t))$ are defined on an extension of $\mathcal{B}$ and mutually independent as well as independent of $\mathcal{F}$. On the other hand, provided that $\Upsilon \equiv 0$ (so the noise is absent), the corresponding result to our estimator can be written as follows:

$$
\begin{align*}
& n^{1 / 4}\left(\operatorname{MRC}[Y]_{t}^{n}-[X, X]_{t}\right) \\
& \quad \rightarrow d_{s} \frac{\sqrt{2 \Phi_{22} \theta}}{\psi_{2}}\left(\sqrt{2} \int_{0}^{t} \sigma_{s}^{2} \sqrt{G_{s}} \mathrm{~d} W_{s}^{\prime}+2 \sum_{r: S_{r} \leq t} \Delta X_{S_{r}} \sigma_{S_{r}} \sqrt{G_{S_{r}}} \Psi_{r}\right), \tag{3.6}
\end{align*}
$$

where we also assume that $G$ is continuous for simplicity. Compared with the above equation with (3.5), the quantities $H^{\prime}$ and $\eta$ coming from the irregularity of the observation times in the left-hand side of (3.5) are replaced with $G$ in (3.6). Since the quantity $H$ contains the information of the second moments of the durations and $\eta$ contains that of all the moments of the durations around the jump times, the distributional future of the durations strongly affects the asymptotic distribution in (3.5). In contrast, the first moments of the durations only affect the asymptotic distribution in (3.6).

Remark 3.11 (Comparison with the continuous case). The result of the theorem is not new if $X$ is continuous. In fact, in the case that $X$ is continuous, a central limit theorem for the MRC estimator can be derived with a somewhat weaker assumption on the limiting process $G$; see Theorem 3.1 and assumption [A4] of Koike [36] for details. In the discontinuous case, we need some regularity of the path of the left limit process $G_{-}$to verify the approximation given in Proposition 6.6, so the structural assumption [A4](iii-c) is necessary.

It is worth mentioning that the structural assumption on $G$ is necessary to deal with the irregularity of observation times in the discontinuous case. In contrast, such a condition is only required to handle the time endogeneity in the continuous case. In fact, if the observation times have a kind of pre-determination property (the so-called strong predictability), convergence in probability of $G^{n}$ to $G$ for the Skorokhod topology is sufficient to derive a central limit theorem; see [35] for details.

Remark 3.12 (Feasible limit theorem). In order to apply Theorem 3.1 to real statistical problems such as the construction of confidence intervals, we need an estimator for the asymptotic covariance matrix given by (3.3) and (3.4). This will be achieved by combining the technique used in the non-synchronously observed diffusion setting (e.g., a kernel approach of [24] or a histogram-type method of [8]) with the one used in the jump diffusion setting (e.g., a thresholding and locally averaging method of Aït-Sahalia, Jacod and Li [2]). Or we can presumably use an estimator of [3] for the equidistant sampling setting without modification because the distribution of the variable $n^{1 / 4} \bar{Y}_{i}$ is, roughly speaking, approximated by the $d$-dimensional normal variable with mean 0 and covariance matrix $\theta \psi_{2} \Sigma_{T_{i}} G_{T_{i}}+\frac{\psi_{1}}{\theta} \widetilde{\Upsilon}_{T_{i}}$ in the absence of jumps conditionally on $\mathcal{F}_{T_{i}}^{(0)}$, where $\widetilde{\Upsilon}$ is the same one as in Remark 3.8 (this is theoretically manifested by Lemma 6.7 in a sense).

## 4. Examples of the observation times

In this section, we give some illustrative examples of the observation times that satisfy the condition [A1]. We shall start to discuss univariate examples (i.e., we assume that $d=1$ ), which are not encompassed with the restricted discretization schemes.

Example 4.1. As an illustrative example of endogenous observation times, we consider a simple model generated by hitting times of the underlying Brownian motion $W$. This type of model is commonly used in the literature; see $[18,37,48]$ among others. Here, we especially focus on a simpler version of the specification from [48]. Specifically, $t_{i}^{1}$ 's are defined as follows:

$$
t_{0}^{1}=0, \quad t_{i+1}^{1}=\inf \left\{t>t_{i}^{1}: W_{t}-W_{t_{i}^{1}}+\sqrt{n} \mathfrak{a}_{t_{i}^{1}}\left(t-t_{i}^{1}\right)=\mathfrak{b}_{t_{i}^{1}} / \sqrt{n}\right\},
$$

where $\mathfrak{a}$ and $\mathfrak{b}$ are cadlag $\left(\mathcal{F}_{t}^{(0)}\right)$-adapted processes such that $\mathfrak{a}_{t} \mathfrak{b}_{t}>0$ and $\mathfrak{a}_{t-} \mathfrak{b}_{t-}>0$ for every $t$. In this case, [A1] is satisfied with setting $T_{p}=\tau_{p}^{1}=t_{p}^{1}$ for every $p$, as long as $G:=\mathfrak{b} / \mathfrak{a}$ satisfies [A1](iii-c). In fact, noting that, conditionally on $\mathcal{F}_{T_{p}}^{(0)}, n\left(T_{p+1}-T_{p}\right)$ follows the inverse Gaussian distribution with mean $G_{T_{p}}$ and variance $G_{T_{p}}^{2} / \mathfrak{a}_{T_{p}}$, (3.1) holds true for any $t>0$ and $\xi \in(0,1)$. Moreover, we have $E\left[n\left(T_{p+1}-T_{p}\right) \mid \mathcal{G}_{T_{p}}^{(0)}\right]=G_{T_{p}}$ for every $p$ because $W$ is independent of $\mu$. Hence, [A1](i)-(iii) are satisfied with $\mathcal{N}^{n}=\varnothing$. Finally, [A1](iv) is automatically satisfied.

Example 4.2. We can also accommodate observation times generated by hitting times of a Brownian motion plus finitely many jumps to our situation. For example, let us consider the observation times defined as follows:

$$
t_{0}^{1}=0, \quad t_{i+1}^{1}=\inf \left\{t>t_{i}^{1}: W_{t}-W_{t_{i}^{1}}+\sqrt{n} \mathfrak{a}_{t_{i}^{1}}\left(t-t_{i}^{1}\right)+\delta^{\prime} \star \mu_{t}-\delta^{\prime} \star \mu_{t_{i}^{1}}=\mathfrak{b}_{t_{i}^{1}} / \sqrt{n}\right\},
$$

where $\mathfrak{a}$ and $\mathfrak{b}$ are the same one as in Example 4.1 and $\delta^{\prime}$ is an $\left(\mathcal{F}_{t}^{(0)}\right)$-optional real-valued function on $\Omega^{(0)} \times \mathbb{R}_{+} \times E$ such that $1_{\left\{\delta^{\prime} \neq 0\right\}} \star \mu_{t}<\infty$ for all $t$. Therefore, the process $\delta^{\prime} \star \mu$ has finitely many jumps. Then, it can easily been seen that [A1] is satisfied in this case under the same situation as that of Example 4.1, except for setting $\mathcal{N}^{n}=\left\{p \in \mathbb{Z}_{+}: \delta^{\prime} \star \mu_{T_{p+1}}-\delta^{\prime} \star \mu_{T_{p}}>0\right\}$.

Example 4.3. Let us consider the observation times discussed in Example 3 of [12]. Namely, $t_{i}^{1}=i / n$ if $i$ is even and $t_{i}^{1}=(i+\alpha) / n$ if $i$ is odd, where $\alpha \in(0,1)$ is a constant. [12] showed that this observation times produce an additional randomness in the asymptotic distribution of the realized covariance estimator even though they are deterministic. In fact, in this case the variable $\eta(t)$ in (3.5) takes the values $(1+\alpha)$ and $(1-\alpha)$ with probabilities $(1+\alpha) / 2$ and $(1-\alpha) / 2$, respectively. On the other hand, setting $T_{p}=(p+1) / n$ and $\tau_{p}^{1}=t_{p}^{1}$, [A1] is satisfied. Hence, in our case this example has the same impact as that of the regular observation times on the asymptotic distribution.

Next, we turn to the multivariate and non-synchronous examples. As the data synchronization method, we focus on the refresh sampling method.

Example 4.4. We shall discuss the Poisson sampling, which is one of the most popular models in this area; see, for example, $[8,12,23,53]$. Let $\left(t_{i}^{k}\right)$ be a sequence of Poisson arrival times with the intensity $n p_{k}$ for each $k$ and suppose that $\left(t_{i}^{1}\right), \ldots,\left(t_{i}^{d}\right)$ are mutually independent and independent of $X$ and $\epsilon$. Then, [A1] is satisfied with

$$
G_{s} \equiv \sum_{k=1}^{d} \sum_{1 \leq l_{1}<\cdots<l_{k} \leq d} \frac{(-1)^{k-1}}{p_{l_{1}}+\cdots+p_{l_{k}}}, \quad \chi_{s}^{k l} \equiv \begin{cases}1, & \text { if } k=l \\ 0, & \text { otherwise }\end{cases}
$$

Example 4.5. Here, we give an example of observation times which are possibly endogenous and satisfy [A1] with the explicit $G$ and $\chi$. More precisely, we give a continuous time analog of the Lo-MacKinlay model of [39].

Let $\left(\tau_{i}\right)_{i=0}^{\infty}$ be a sampling scheme and suppose that $\sup _{i \geq 0}\left(\tau_{i} \wedge t-\tau_{i-1} \wedge t\right)=\mathrm{o}_{p}\left(n^{-\xi}\right)$ as $n \rightarrow \infty$ for any $t>0$ and $\xi \in(0,1)$. For each $k=1, \ldots, d$, let $\left(M_{k}(n, i)\right)_{i=0}^{\infty}$ be a sequence of $\mathbb{Z}_{+}$-valued variables defined on $\mathcal{B}^{(0)}$ and independent of $X$ and $\left(\tau_{i}\right)$ such that $\mathcal{M}_{k}:=\left(M_{k}(n, i+1)-M_{k}(n, i)\right)_{i=0}^{\infty}$ is independent and geometrically distributed with the common success probability $p_{k} \in(0,1)$. Moreover, suppose that, for each $i, M_{k}(n, i)$ is an $\left(\mathcal{F}_{\tau_{j}}^{(0)}\right)_{j=0^{-}}^{\infty}$ stopping time so that $t_{i}^{k}:=\tau_{M_{k}(n, i)}$ is an $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping time and that $M_{k}(n, i+1)-M_{k}(n, i)$ is independent of $\mathcal{F}_{t_{i}^{k}}^{(0)}$. Finally, assume that $\mathcal{M}_{1}, \ldots, \mathcal{M}_{d}$ are mutually independent and that [A1](i)-(iv) are satisfied with replacing $\left(T_{i}\right)$ by $\left(\tau_{i}\right)$. Then it can easily be shown that [A1] holds true with

$$
G_{s} \equiv \sum_{k=1}^{d} \sum_{1 \leq l_{1}<\cdots<l_{k} \leq d} \frac{(-1)^{k-1} G_{s}^{0}}{1-\left(1-p_{l_{1}}\right) \cdots\left(1-p_{l_{k}}\right)}
$$

and

$$
\chi_{s}^{k l} \equiv \begin{cases}1, & \text { if } k=l \\ p_{k} p_{l} /\left(p_{k}+p_{l}-p_{k} p_{l}\right), & \text { otherwise }\end{cases}
$$

Here, $G^{0}$ denotes the asymptotic conditional expected duration process corresponding to ( $\tau_{i}$ ). By taking an endogenous sampling scheme as the underlying sampling scheme ( $\tau_{i}$ ), we can obtain endogenous observation times.

## 5. Simulation study

In this section, we assess the finite sample accuracy of the central limit theory developed in this paper and confirm our theoretical findings via Monte Carlo experiments.

We simulate over the unit interval [0, 1], and basically follow the design of [3]. To simulate the latent semimartingale $X$, the following bivariate Heston model with jumps is considered:

$$
\mathrm{d} X_{t}^{k}=\sigma_{k, t} \mathrm{~d} W_{t}^{k}+\mathrm{d} Z_{t}^{k}, \quad \mathrm{~d} \sigma_{k, t}^{2}=\kappa_{k}\left(\bar{\sigma}_{k}^{2}-\sigma_{k, t}^{2}\right) \mathrm{d} t+s_{k} \sigma_{k, t} \mathrm{~d} B_{t}^{k}+\mathrm{d} J_{t}^{k}-\lambda_{k}^{V} \tau_{k}^{V} \mathrm{~d} t
$$

$$
k=1,2 .
$$

Here, $W^{1}, W^{2}, B^{1}, B^{2}$ are correlated standard Brownian motions such that

$$
\begin{aligned}
& \mathrm{d}\left[W^{1}, W^{2}\right]_{t}=\rho_{B} \mathrm{~d} t, \quad \mathrm{~d}\left[W^{k}, B^{k}\right]_{t}=\rho_{k} \mathrm{~d} t, \\
& \mathrm{~d}\left[W^{1}, B^{2}\right]_{t}=\mathrm{d}\left[W^{2}, B^{1}\right]_{t}=\mathrm{d}\left[B^{1}, B^{2}\right]_{t}=0 .
\end{aligned}
$$

$J^{k}$ is a compound Poisson process with jump size uniformly distributed on $\left[0,2 \tau_{k}^{V}\right]$ and jump intensity $\lambda_{k}^{V} . J^{1}$ and $J^{2}$ are assumed to be mutually independent. $Z^{k}$ is a pure jump Lévy process specified as follows. First, $Z^{2}$ is linearly correlated with $Z^{1}$ as $Z^{2}=\rho_{J} Z^{1}+\sqrt{1-\rho_{J}^{2}} Z^{0}$, where $Z^{0}$ is another Lévy process independent of $Z^{1}$. For each $m=0,1, Z^{m}$ is a CGMY process with Lévy density given by

$$
f_{m}(x)=c_{m} \frac{\mathrm{e}^{-\gamma_{m-}-x \mid}}{|x|^{1+\beta_{m}}} 1_{\{x<0\}}+c_{m} \frac{\mathrm{e}^{-\gamma_{m+} x}}{x^{1+\beta_{m}}} 1_{\{x>0\}} .
$$

The parameter values of the stochastic volatility processes used in the simulation are reported in Table 1. The initial value for the volatility processes $\sigma_{k, t}^{2}$ is set at $\bar{\sigma}_{k}^{2}$ for each $k=1,2$, which ensures that $E\left[\sigma_{k, t}^{2}\right]=\bar{\sigma}_{k}^{2}$ for all $t \in[0,1]$. The specification of the parameters in the CGMY processes is as follows. We set $\gamma_{m+}=3, \gamma_{m-}=5, \beta_{m}=0.5$ for every $m=0,1 . c_{1}$ is selected such that the quadratic variation contributed by jumps in $X^{1}$ amounts to $15 \%$ in expectation, that is, $E\left(\left[Z^{1}, Z^{1}\right]_{1}\right) / E\left(\left[X^{1}, X^{1}\right]_{1}\right)=0.15$. Then $c_{0}$ is selected such that $E\left(\left[Z^{2}, Z^{2}\right]_{1}\right) / E\left(\left[X^{2}, X^{2}\right]_{1}\right)=0.15$. Finally, the correlation parameter $\rho_{J}$ between the jump processes are set at 0.2 . Note that $Z^{1}$ and $Z^{2}$ can be exactly simulated because we only consider the situation where they are of finite variation; see, for example, [33] for details.

To generate observation times, we consider Lo-MacKinlay type sampling schemes illustrated in Example 4.5. Two kinds of sequence $\left(\tau_{i}\right)_{i=0}^{\infty}$ of latent observation times are considered: one is the equidistant sampling scheme $\tau_{i}=i / n$ and the other is the endogenous sampling scheme defined by

$$
\begin{equation*}
\tau_{0}=0, \quad \tau_{i+1}=\inf \left\{t>\tau_{i}: W_{t}^{1}-W_{\tau_{i}}^{1}-2 \sqrt{n}\left(t-\tau_{i}\right)=-2 / \sqrt{n}\right\}, \quad i=0,1, \ldots, \tag{5.1}
\end{equation*}
$$

where we set $n=23400$. Note that in the latter case the sequence $\left(\tau_{i+1}-\tau_{i}\right)_{i=0}^{\infty}$ is independent and identically distributed with the inverse Gaussian distribution with mean $1 / n$ and variance $4 / n^{2}$, thus we can exactly simulate $\tau_{i}$ 's (and construct the exactly discretized path $\left\{W_{\tau_{i}}\right\}$ from $\left\{\tau_{i}\right\}$ ). Furthermore, in both cases the corresponding conditional expected duration processes $G^{0}$ are identical with 1 . The parameters $p_{1}$ and $p_{2}$ from Example 4.5, which denote the proba-

Table 1. The parameters of the stochastic volatility processes

| $k$ | $\kappa_{k}$ | $s_{k}$ | $\bar{\sigma}_{k}$ | $\rho_{k}$ | $\lambda_{k}^{V}$ | $\tau_{k}^{V}$ | $\rho_{B}$ |
| :--- | :--- | :--- | :--- | :--- | ---: | :--- | :--- |
| 1 | 5 | 0.3 | 0.25 | -0.6 | 5 | 0.05 | 0.5 |
| 2 | 4 | 0.4 | 0.3 | -0.75 | 10 | 0.01 | - |

bilities of observations occurring, are assumed to be identical each other and varied thorough $1 / 3,1 / 5,1 / 10$ and $1 / 30$.

In constructing noisy prices $Y$, we first generate a discretized path $X_{\tau_{0}}, X_{\tau_{1}}, \ldots$ of $X$ using a standard Euler scheme. After that, we add simulated microstructure noise $Y_{\tau_{i}}=X_{\tau_{i}}+\epsilon_{\tau_{i}}$ by generating centered Gaussian i.i.d. variables $\epsilon_{\tau_{0}}^{k}, \epsilon_{\tau_{1}}^{k}, \ldots$ with standard deviation 0.005. $\epsilon^{1}$ and $\epsilon^{2}$ are assumed to be mutually independent. Simulation results are based on 10000 Monte Carlo iterations for each scenario.

Following [14], the MRC estimator is implemented using the weight function $g(x)=x \wedge(1-$ $x$ ) and the refresh time sampling method (the finite sample corrections explained in [14] are also included). We consider the window size $k_{n}$ of the form $k_{n}=\left\lceil\theta \sqrt{N_{1}^{n}}\right\rceil$, and $\theta$ is selected among $1 / 3$ and 1 . The former value of $\theta$ corresponds to the one used in [28], while the latter one does to the one used in [14]. We assess the accuracy of the standard normal approximation of the infeasible standardized statistic

$$
\begin{equation*}
n^{1 / 4} \frac{\mathrm{MRC}[Y]_{1}^{n, 12}-\left[X^{1}, X^{2}\right]_{1}}{\sqrt{\text { AVAR }}} \tag{5.2}
\end{equation*}
$$

where AVAR is the theoretical asymptotic variance given in Theorem 3.1. Table 2 reports the sample mean and standard deviation as well as $95 \%$ and $99 \%$ coverages of (5.2). As the table reveals, the central limit theorem for (5.2) fairly works. As was expected from the theory developed in the above, we find no significant difference of the results between the exogenous and the endogenous sampling cases. At relatively low frequencies like $p_{1}=p_{2}=1 / 10$ or $1 / 30$, the results for $\theta=1$ show the better performance than those for $\theta=1 / 3$. This would be because $k_{n}$ is not sufficiently large in such a situation, in order to work the averaging effect of the pre-averaging procedure explained in Remark 3.7.

Table 2. Simulation results of the standardized estimates

|  | $\tau_{i}=i / n$ |  |  |  | $\tau_{i}$ 's are defined by (5.1) |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Mean | SD | $\begin{aligned} & \text { Coverage } \\ & (95 \%) \end{aligned}$ (95\%) | Coverage (99\%) | Mean | SD | Coverage (95\%) | Coverage (99\%) |
| $\theta=1 / 3$ |  |  |  |  |  |  |  |  |
| $p_{1}=p_{2}=1 / 3$ | -0.00 | 1.01 | 0.949 | 0.987 | -0.00 | 1.02 | 0.948 | 0.987 |
| $p_{1}=p_{2}=1 / 5$ | -0.01 | 1.02 | 0.946 | 0.987 | -0.01 | 1.04 | 0.943 | 0.986 |
| $p_{1}=p_{2}=1 / 10$ | -0.01 | 1.05 | 0.939 | 0.983 | -0.01 | 1.06 | 0.937 | 0.984 |
| $p_{1}=p_{2}=1 / 30$ | -0.03 | 1.09 | 0.928 | 0.979 | -0.03 | 1.10 | 0.929 | 0.980 |
| $\theta=1$ |  |  |  |  |  |  |  |  |
| $p_{1}=p_{2}=1 / 3$ | -0.01 | 1.01 | 0.948 | 0.987 | -0.01 | 1.01 | 0.949 | 0.989 |
| $p_{1}=p_{2}=1 / 5$ | -0.01 | 1.01 | 0.948 | 0.987 | -0.01 | 1.01 | 0.951 | 0.987 |
| $p_{1}=p_{2}=1 / 10$ | -0.02 | 1.02 | 0.947 | 0.986 | -0.02 | 1.02 | 0.948 | 0.987 |
| $p_{1}=p_{2}=1 / 30$ | -0.03 | 1.03 | 0.946 | 0.985 | -0.03 | 1.03 | 0.943 | 0.985 |

[^0]
## 6. Proof of Theorem 3.1

### 6.1. Preliminaries

### 6.1.1. Localization

Before starting the proof, we strengthen our assumptions [A1]-[A3] by localization procedures. First, a standard localization procedure, described in detail in Lemma 4.4.9 of [31], for instance, allows us to replace the conditions [A2] and [A3] by the following strengthened versions, respectively:
[SA2] We have [A2], and the processes $X_{t}, b_{t}, \sigma_{t}, \widetilde{b}_{t}$ and $\widetilde{\sigma}_{t}$ are bounded. Also, $b_{t}$ is $\left(\mathcal{H}_{t}^{\wedge}\right)$-progressively measurable and $\sigma_{t}$ is $\left(\mathcal{H}_{t}^{\wedge}\right)$-adapted. Moreover, there are a constant $\Lambda$ and a non-negative bounded function $\gamma$ on $E$ such that $\int \gamma(z)^{2} \lambda(\mathrm{~d} z)<\infty$ and $\left\|\delta\left(\omega^{(0)}, t, z\right)\right\| \vee$ $\left\|\widetilde{\delta}\left(\omega^{(0)}, t, z\right)\right\| \leq \gamma(z)$ and

$$
\begin{aligned}
E\left[\left\|b_{t_{1}}-b_{t_{2}}\right\|^{2} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right] & \leq \Lambda E\left[\left|t_{1}-t_{2}\right|^{\sigma} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right], \\
E\left[\left\|\delta\left(t_{1}, z\right)-\delta\left(t_{2}, z\right)\right\|^{2} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right] & \leq \Lambda \gamma(z)^{2} E\left[\left|t_{1}-t_{2}\right|^{\Phi} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right]
\end{aligned}
$$

for any bounded $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping times $t_{1}$ and $t_{2}$.
[SA3] There are a constant $\Gamma>4$ and a constant $\Lambda^{\prime}$ such that the process $\int\|z\|^{\Gamma} Q_{t}(\mathrm{~d} z)$ is bounded and

$$
E\left[\left\|\Upsilon_{t_{1}}-\Upsilon_{t_{2}}\right\|^{2} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right] \leq \Lambda^{\prime} E\left[\left|t_{1}-t_{2}\right|^{\sigma} \mid \mathcal{F}_{t_{1} \wedge t_{2}}\right]
$$

for any bounded $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping times $t_{1}$ and $t_{2}$. Moreover, $\Upsilon_{t}$ is cadlag and $\left(\mathcal{H}_{t}^{\wedge}\right)$-adapted.
Next, we introduce a strengthened version of [A1]. In the following, we fix a constant $\xi \in$ $(0,1)$ such that

$$
\begin{equation*}
\xi>\frac{7}{8} \vee \frac{1}{2}\left(\kappa+\frac{3}{2}\right) \vee(1-\varpi), \tag{6.1}
\end{equation*}
$$

and we set $\bar{r}_{n}=n^{-\xi}$.
[SA1] We have [A1], and for every $n$ it holds that

$$
\begin{equation*}
\sup _{p \geq 0}\left(T_{p}-T_{p-1}\right) \leq \bar{r}_{n} . \tag{6.2}
\end{equation*}
$$

The following lemma allows us to replace [A1] by [SA1] via another localization argument. The proof is similar to that of Lemma 6.3 from [36], so we omit it.

Lemma 6.1. Assume [A1]. One can find sampling schemes $\left(\widetilde{T}_{p}\right)$ and $\left(\widetilde{\tau}_{p}^{k}\right)(k=1, \ldots, d)$ satisfying the following conditions:
(i) $\left(\widetilde{T}_{p}\right)$ and $\left(\widetilde{\tau}_{p}^{k}\right)$ satisfy [SA1] with the same limiting processes $G$ and $\chi$ as those of the original sampling schemes.
(ii) For any $t>0$ there is a subset $\Omega_{n, t}^{(0)}$ of $\Omega^{(0)}$ such that $\lim _{n} P^{(0)}\left(\Omega_{n, t}^{(0)}\right)=1$. Moreover, on $\Omega_{n, t}^{(0)}$ we have $T_{p} \wedge t=\widetilde{T}_{p} \wedge t$ and $\tau_{p}^{k} \wedge t=\widetilde{\tau}_{p}^{k} \wedge t$ for all $k, p$.

### 6.1.2. Outline of the proof

Here, we give a brief description of the scheme of the proof. First, for the proof it is convenient to realize the processes $\mathcal{W}$ and $\mathcal{Z}$ on an extension of $\mathcal{B}^{(0)}$ as in Remarks 3.8-3.9 (so we will use the notation introduced in these remarks in the following). For notational simplicity, we use the same letters $P$ and $E$ for the probability and the expectation with respect to this extension.

Next we introduce some notation. We denote by $\mathcal{R}_{m}$ the set of all indices $r$ such that $S_{r}=$ $S\left(m^{\prime}, j\right)$ for some $j \geq 1$ and some $m^{\prime} \leq m$. Also, we set

$$
\left\{\begin{array}{l}
b(m)_{t}=b_{t}-\int_{A_{m} \cap\{z:|\delta(t, z)| \leq 1\}} \delta(t, z) \lambda(\mathrm{d} z), \quad B(m)_{t}=\int_{0}^{t} b(m)_{s} \mathrm{~d} s, \\
M_{t}=\int_{0}^{t} \sigma_{s} \mathrm{~d} W_{s}, \\
C(m)_{t}=X_{0}+B(m)_{t}+M_{t}, \quad J(m)_{t}=\delta 1_{A_{m}} \star \mu_{t}, \quad X(m)_{t}=C(m)_{t}+J(m)_{t}, \\
Z(m)_{t}=X_{t}-X(m)_{t}=\delta 1_{A_{m}^{c}} \star(\mu-v)_{t} .
\end{array}\right.
$$

These processes are well-defined under [SA2]. Furthermore, set $I_{p}=\left[T_{p-1}, T_{p}\right.$ ) for every $p \in$ $\mathbb{Z}_{+}$. On the other hand, for any process $V$ and any (random) interval $I=[S, T)$, we define the random variable $V(I)$ by $V(I)=V_{T}-V_{S}$. We also set $I(t)=I \cap[0, t)=[S \wedge t, T \wedge t)$ for any $t \in \mathbb{R}_{+}$and $|I|=T-S$. For any real-valued function $u$ on $[0,1]$, we set $u_{p}^{n}=u\left(p / k_{n}\right)$ for $p=$ $0,1, \ldots, k_{n}$. For any $d$-dimensional processes $U, V$, any $k, l \in\{1, \ldots, d\}$ and any $u, v \in\left\{g, g^{\prime}\right\}$, we define the process $\Xi_{u, v}^{(k, l)}(U, V)^{n}$ by

$$
\Xi_{u, v}^{(k, l)}(U, V)_{t}^{n}=\frac{1}{\psi_{2} k_{n}} \sum_{i=1}^{N_{t}^{n}-k_{n}+1} \bar{U}(u)_{i}^{k} \bar{V}(v)_{i}^{l}, \quad t \in \mathbb{R}_{+}
$$

where $\bar{U}(u)_{i}^{k}=\sum_{p=0}^{k_{n}-1} u_{p}^{n} U^{k}\left(I_{i+p}\right)$ and $\bar{V}(v)_{i}^{l}$ is defined analogously. Moreover, we define the $d$-dimensional process $\mathfrak{E}$ by

$$
\mathfrak{E}_{t}^{k}=-\frac{1}{k_{n}} \sum_{p=1}^{\infty} \epsilon_{\tau_{p}^{k}}^{k} 1_{\left\{\tau_{p}^{k} \leq t\right\}}, \quad t \in \mathbb{R}_{+}, k=1, \ldots, d
$$

It can easily been seen that $\mathfrak{E}$ is a purely discontinuous locally square-integrable martingale on $\mathcal{B}$ under [SA3]. Finally, for any $d$-dimensional process $V$ we define the $\mathbb{R}^{d} \otimes \mathbb{R}^{d}$-valued process $\Xi[V]^{n}$ by

$$
\Xi[V]^{n, k l}=\Xi_{g, g}^{(k, l)}(V, V)^{n}+\Xi_{g, g^{\prime}}^{(k, l)}(V, \mathfrak{E})^{n}+\Xi_{g, g^{\prime}}^{(l, k)}(V, \mathfrak{E})^{n}+\Xi_{g^{\prime}, g^{\prime}}^{(k, l)}(\mathfrak{E}, \mathfrak{E})^{n}, \quad k, l=1, \ldots, d
$$

Now we turn to the outline of the proof. In the first step, we show that the errors from the interpolations to the synchronized sampling times are asymptotically negligible:

Proposition 6.1. Assume [SA1]-[SA3] and [A4](ii). Then $n^{1 / 4}\left(\mathrm{MRC}[Y]^{n}-\Xi[X]^{n}\right.$ $\left.+\frac{\psi_{1}}{\psi_{2} k_{n}^{2}}[Y, Y]^{n}\right) \xrightarrow{\text { ucp }} 0$.

The proof of this proposition is an easy extension of that of Proposition 6.1 from [36], so we omit it.

In the next step, we decompose the quantity $\boldsymbol{\Xi}[X]^{n}$ as $\boldsymbol{\Xi}[X]^{n}=\boldsymbol{\Xi}[X(m)]_{t}^{n}+\left(\Xi[X]^{n}-\right.$ $\Xi[X(m)]^{n}$ ) for each $m$, and show that the first term enjoys a central limit theorem for any fixed $m$ and the second term is negligible as $m \rightarrow \infty$. More precisely, we prove the following propositions.

Proposition 6.2. Suppose that [SA1]-[SA3] and [A4] are satisfied. Then

$$
n^{1 / 4}\left(\boldsymbol{\Xi}[X(m)]_{t}^{n}-[X(m), X(m)]_{t}-\frac{\psi_{1}}{\psi_{2} k_{n}^{2}}[Y, Y]_{t}^{n}\right) \rightarrow^{d_{s}} \mathcal{W}_{t}+\mathcal{Z}(m)_{t}
$$

as $n \rightarrow \infty$ for any $t>0$ and any $m \geq 1$, where $\mathcal{Z}(m)_{t}=\sum_{r \in \mathcal{R}_{m}: S_{r} \leq t}\left(\mathfrak{Z}_{r}+\mathfrak{Z}_{r}^{*}\right)$.
When further $X$ is continuous, the processes $n^{1 / 4}\left(\Xi[X]^{n}-[X, X]-\frac{\psi_{1}}{\psi_{2} k_{n}^{2}}[Y, Y]^{n}\right)$ converge stably in law to the process $\mathcal{W}$ for the Skorokhod topology.

Proposition 6.3. Suppose that [SA1]-[SA3] and [A4](ii) are satisfied. Then

$$
\begin{equation*}
\mathcal{Z}(m)_{t} \rightarrow^{d_{s}} \mathcal{Z}_{t} \tag{6.3}
\end{equation*}
$$

as $m \rightarrow \infty$ and

$$
\begin{equation*}
\limsup _{m \rightarrow \infty} \limsup _{n \rightarrow \infty} P\left(n^{1 / 4}\left\|\boldsymbol{\Xi}[X]_{t}^{n}-\boldsymbol{\Xi}[X(m)]_{t}^{n}\right\|>\eta\right)=0 \tag{6.4}
\end{equation*}
$$

for any $t, \eta>0$.
Combining Propositions 6.1-6.3 with Proposition 2.2.4 of [31], we obtain Theorem 3.1.

### 6.2. Proof of Proposition 6.2

Throughout the discussions, for (deterministic) sequences $\left(x_{n}\right)$ and ( $y_{n}$ ), $x_{n} \lesssim y_{n}$ means that there is a (non-random) constant $K \in[0, \infty)$ such that $x_{n} \leq K y_{n}$ for large $n$. We also denote by $E_{0}$ the conditional expectation given $\mathcal{F}^{(0)}$, that is, $E_{0}[\cdot]=E\left[\cdot \mid \mathcal{F}^{(0)}\right]$.

The proof of Proposition 6.2 is divided into the following steps:
(i) Approximating the estimation error due to the diffusion part by a more tractable one,
(ii) Proving a central limit theorem for the approximation constructed in (i),
(iii) Approximating the estimation error due to the jump part by a more tractable one (Section 6.2.1),
(iv) Proving a local stable convergence result corresponding to Lemma 16.3.7 of [31] (Section 6.2.2),
(v) Proving a joint limit theorem for the pair of the above approximations and completing the proof of the proposition (Section 6.2.3).

The first two steps have already been carried out in [36]. For the later use, we summarize the result in the following. We begin by introducing some notation. For any $d$-dimensional processes $U, V$, any $k, l \in\{1, \ldots, d\}$ and any real-valued functions $u, v$ on $[0,1]$, we define the processes $\mathbb{M}_{u, v}^{(k, l)}(U, V)^{n}$ and $\mathbb{L}_{u, v}^{(k, l)}(U, V)^{n}$ by

$$
\mathbb{M}_{u, v}^{(k, l)}(U, V)_{t}^{n}=\sum_{q=2}^{N_{t}^{n}+1} C_{u, v}^{n}(U)_{q}^{k} V^{l}\left(I_{q}\right), \quad \mathbb{L}_{u, v}^{(k, l)}(U, V)_{t}^{n}=\mathbb{M}_{u, v}^{(k, l)}(U, V)_{t}^{n}+\mathbb{M}_{v, u}^{(l, k)}(V, U)_{t}^{n},
$$

where

$$
C_{u, v}^{n}(U)_{q}^{k}=\sum_{p=\left(q-k_{n}\right) \vee 1}^{q-1} c_{u, v}^{n}(p, q) U^{k}\left(I_{p}\right), \quad c_{u, v}^{n}(p, q)=\frac{1}{\psi_{2} k_{n}} \sum_{i=\left(p \vee q-k_{n}+1\right) \vee 1}^{p \wedge q} u_{p-i}^{n} v_{q-i}^{n} .
$$

Moreover, define the $\mathbb{R}^{d} \otimes \mathbb{R}^{d}$-valued process $\mathbf{L}[M]^{n}$ by

$$
\mathbf{L}[M]^{n, k l}=\mathbb{L}_{g, g}^{(k, l)}(M, M)^{n}+\mathbb{L}_{g, g^{\prime}}^{(k, l)}(M, \mathfrak{E})^{n}+\mathbb{L}_{g, g^{\prime}}^{(l, k)}(M, \mathfrak{E})^{n}+\mathbb{L}_{g^{\prime}, g^{\prime}}^{(k, l)}(\mathfrak{E}, \mathfrak{E})^{n} .
$$

Then, we have the following results, which are proved as Proposition 6.2 and equation (6.7) from [36]:

Proposition 6.4. Suppose that [SA1]-[SA3] and [A4](ii) are satisfied. Then

$$
n^{1 / 4}\left(\boldsymbol{\Xi}[C(m)]^{n}-\mathbf{L}[M]^{n}-[M, M]-\frac{\psi_{1}}{\psi_{2} k_{n}^{2}}[Y, Y]^{n}\right) \xrightarrow{\text { ucp }} 0
$$

as $n \rightarrow \infty$ for any $m \geq 1$.
Proposition 6.5. Suppose that [SA1]-[SA3] and [A4](ii) are satisfied. Then the processes $n^{1 / 4} \mathbf{L}[M]^{n}$ converge stably in law to $\mathcal{W}$ for the Skorokhod topology.

From the next section, we start the proofs of the remaining steps.

### 6.2.1. Approximation of the estimation error due to the jump part

In this subsection we fix $t>0$ and $m \in \mathbb{N}$, and denote by $\Omega_{n}(t, m)$ the set on which $k_{n}-1 \leq$ $N_{S_{r}-}^{n} \leq N_{t}^{n}-k_{n}$ for all $r \in \mathcal{R}_{m}$ such that $S_{r} \leq t$. On this set, we have

$$
\begin{align*}
n^{1 / 4} \Xi_{g, g}^{(k, l)}(C(m), J(m))_{t}^{n} & =\frac{n^{1 / 4}}{\psi_{2} k_{n}} \sum_{i=0}^{N_{t}^{n}-k_{n}+1} \sum_{p, q=0}^{k_{n}-1} g_{p}^{n} g_{q}^{n} C(m)^{k}\left(I_{i+p}\right) J(m)^{l}\left(I_{i+q}\right)  \tag{6.5}\\
& =\sum_{r \in \mathcal{R}_{m}: S_{r} \leq t}\left\{\eta_{+}(n, r)^{k}+\eta_{-}(n, r)^{k}\right\} \Delta X_{S_{r}}^{l}
\end{align*}
$$

where

$$
\left\{\begin{array}{l}
\eta_{+}(n, r)=n^{1 / 4} \sum_{p=N_{S_{r}-}^{n}}^{N_{S_{r}-}^{n}+k_{n}} c_{g, g}^{n}\left(p, N_{S_{r}-}^{n}+1\right) C(m)\left(I_{p}\right), \\
\eta_{-}(n, r)=n^{1 / 4} \sum_{p=\left(N_{\left.S_{r_{-}}-k_{n}+2\right)_{+}}^{n} c_{g, g}^{n}\left(p, N_{S_{r}-}^{n}+1\right) C(m)\left(I_{p}\right) .\right.}^{N_{S_{r}-}^{n}} .
\end{array}\right.
$$

Similarly, on $\Omega_{n}(t, m)$ we have

$$
\begin{equation*}
n^{1 / 4} \Xi_{g^{\prime}, g}^{(k, l)}(\mathfrak{E}, J(m))_{t}^{n}=\sum_{r \in \mathcal{R}_{m}: S_{r} \leq t}\left\{\eta_{+}^{\prime}(n, r)^{k}+\eta_{-}^{\prime}(n, r)^{k}\right\} \Delta X_{S_{r}}^{l}, \tag{6.6}
\end{equation*}
$$

where

$$
\left\{\begin{array}{l}
\eta_{+}^{\prime}(n, r)^{k}=-\frac{n^{1 / 4}}{k_{n}} \sum_{p=N_{S_{r}-}^{n}+1}^{N_{S_{r}-}^{n}+k_{n}} c_{g^{\prime}, g}^{n}\left(p, N_{S_{r}-}^{n}+1\right) \epsilon_{\tau_{p}^{k}}^{k}, \\
\eta_{-}^{\prime}(n, r)^{k}=-\frac{n^{1 / 4}}{k_{n}} \sum_{p=\left(N_{S_{r}-}^{n}-k_{n}+2\right)_{+}}^{N_{S_{r}-}^{n}} c_{g^{\prime}, g}^{n}\left(p, N_{S_{r}-}^{n}+1\right) \epsilon_{\tau_{p}^{k}}^{k} .
\end{array}\right.
$$

The aim of this subsection is to approximate $\eta_{ \pm}(n, r)$ and $\eta_{ \pm}^{\prime}(n, r)$ by more tractable quantities. Here, the major difficulty coming from the irregularity of the observation times is the fact that $N_{S_{r}-}^{n}-k_{n}+1$ might not be a $\left(\mathcal{G}_{T_{p}}^{(0)}\right)_{p=0}^{\infty}$-stopping time. Therefore, we first "approximate" $N_{S_{r}-}^{n}-$ $k_{n}+1$ by a $\left(\mathcal{G}_{T_{p}}^{(0)}\right)_{p=0}^{\infty}$-stopping time.

More precisely, set $\underline{S}_{r}=\left(S_{r}-\frac{k_{n}}{n} \log n\right)_{+}$and $S_{r}^{\dagger}=\left(S_{r}-\frac{k_{n}}{n} G_{\underline{S}_{r}}^{n} \wedge \log n\right)_{+}$. Then, $S_{r}^{\dagger}$ is a $\left(\mathcal{G}_{t}^{(0)}\right)$-stopping time (Lemma 6.2), and thus $N_{S_{r}^{\dagger}}^{n}+1$ is a $\left(\mathcal{G}_{T_{p}}^{(0)}\right)_{p=0}^{\infty}$-stopping time, and this variable gives an approximation of $N_{S_{r}-}^{n}-k_{n}+1$ (Lemma 6.4).

Now we can define our tractable approximations of $\eta_{ \pm}(n, r)$ and $\eta_{ \pm}^{\prime}(n, r)$ as follows. For any non-negative random variable $\rho$ and any real-valued function $\phi$ on $[0,1]$, we define the $d^{\prime}$ dimensional variable $L(\phi, \rho)_{n}=\left(L(\phi, \rho)_{n}^{j}\right)_{1 \leq j \leq d^{\prime}}$ and the $d$-dimensional variable $L^{\prime}(\phi, \rho)_{n}=$ $\left(L^{\prime}(\phi, \rho)_{n}^{k}\right)_{1 \leq k \leq d}$ by

$$
L(\phi, \rho)_{n}^{j}=n^{1 / 4} \sum_{w=1}^{k_{n}-1} \phi_{w}^{n} W^{j}\left(I_{i(\rho)^{n}+w}\right), \quad L^{\prime}(\phi, \rho)_{n}^{k}=\frac{n^{1 / 4}}{k_{n}} \sum_{w=1}^{k_{n}-1} \phi_{w}^{n} \epsilon_{\tau_{i(\rho)^{n}+w}^{k}}^{k},
$$

where we set $i(\rho)^{n}=N_{\rho}^{n}+1$ (recall that $\phi_{w}^{n}=\phi\left(w / k_{n}\right)$ ). We also define the function $\widetilde{\phi}$ on $[0,1]$ by $\widetilde{\phi}(x)=\phi(1-x)$. Then we set

$$
\begin{cases}z_{r-}^{n}=\psi_{2}^{-1} L\left(\widetilde{\phi_{g, g}}, S_{r}^{\dagger}\right)_{n}, & z_{r+}^{n}=\psi_{2}^{-1} L\left(\phi_{g, g}, S_{r}\right)_{n}, \\ z_{r-}^{\prime n}=-\psi_{2}^{-1} L^{\prime}\left(\widetilde{\phi_{g^{\prime}, g}}, S_{r}^{\dagger}\right)_{n}, & z_{r+}^{\prime n}=-\psi_{2}^{-1} L^{\prime}\left(\phi_{g^{\prime}, g}, S_{r}\right)_{n}\end{cases}
$$

The aim of this subsection is to prove the following proposition.
Proposition 6.6. Suppose that [SA1]-[SA3] and [A4](ii) are satisfied. Then

$$
\begin{align*}
& \eta_{-}(n, r)=\sigma_{S_{r}^{\dagger}} z_{r-}^{n}+\mathrm{o}_{p}(1), \quad \eta_{+}(n, r)=\sigma_{S_{r}} z_{r+}^{n}+\mathrm{o}_{p}(1),  \tag{6.7}\\
& \eta_{-}^{\prime}(n, r)=z_{r-}^{\prime n}+\mathrm{o}_{p}(1), \quad \eta_{+}^{\prime}(n, r)=z_{r+}^{\prime n}+\mathrm{o}_{p}(1) . \tag{6.8}
\end{align*}
$$

Now we start to justify that the variable $N_{S_{r}^{\dagger}}^{n}+1$ is an appropriate approximation of $N_{S_{r}-}^{n}-$ $k_{n}+1$. In the remainder of this subsection we fix an index $r \in \mathcal{R}_{m}$ such that $S_{r}<\infty$.

Lemma 6.2. Under [SA1], $S_{r}^{\dagger}$ is a $\left(\mathcal{G}_{t}^{(0)}\right)$-stopping time.
Proof. For any $t \geq 0$, we have $\left\{S_{r}^{\dagger} \leq t\right\}=\left\{S_{r} \leq t+\left(k_{n} / n\right) G_{\underline{S}_{r}}^{n} \wedge \log n\right\} \cap\left\{\underline{S}_{r} \leq t\right\}$. Therefore, noting that $S_{r}$ is $\mathcal{G}_{0}^{(0)}$-measurable, we obtain $\left\{S_{r}^{\dagger} \leq t\right\} \in \mathcal{G}_{t}^{(0)}$.

Lemma 6.3. Under [SA1], $\sup _{0<h<h_{0}}\left|G_{S_{r}-}-G_{\left(S_{r}-h\right)_{+}}\right|=\mathrm{O}_{p}\left(\sqrt{h_{0}}\right)$ as $h_{0} \downarrow 0$.
Proof. Define the processes $G(m), G^{\prime}(m)$ and $G^{\prime \prime}(m)$ by $G(m)_{t}=\int_{0}^{t} \widehat{\sigma}_{s} \mathrm{~d} W_{s}+\left(\widehat{\delta} 1_{A_{m}^{c} \cap\{\{\tilde{\delta} \mid \leq 1\}}\right) \star$ $(\mu-\nu)_{t}, G^{\prime}(m)_{t}=\left(\widehat{\delta} 1_{A_{m} \cap\{\{\widetilde{\delta} \mid \leq 1\}}+\widehat{\delta} 1_{\{|\widetilde{\delta}|>1\}}\right) \star \mu_{t}$, and $G^{\prime \prime}(m)=G-G(m)-G^{\prime}(m)$. Since $G^{\prime}(m)$ is piecewise constant, it is evident that $\sup _{0<h<h_{0}}\left|G^{\prime}(m)_{S_{r}-}-G^{\prime}(m)_{\left(S_{r}-h\right)_{+}}\right|=$ $\mathrm{O}_{p}\left(\sqrt{h_{0}}\right)$. Moreover, since $G^{\prime \prime}(m)$ is absolutely continuous with a locally bounded derivative, it also holds that $\sup _{0<h<h_{0}}\left|G^{\prime \prime}(m)_{S_{r}-}-G^{\prime \prime}(m)_{\left(S_{r}-h\right)_{+}}\right|=\mathrm{O}_{p}\left(\sqrt{h_{0}}\right)$. On the other hand, let $\left(\mathcal{G}_{t}^{A_{m}}\right)$ be the smallest filtration containing $\left(\mathcal{F}_{t}^{(0)}\right)$ such that $\mathcal{G}_{0}^{A_{m}}$ contains the $\sigma$-field generated by the restriction of the measure $\mu$ to $\mathbb{R}_{+} \times A_{m}$. Then, by Proposition 2.1.10 of [31] $G(m)$ is a locally square integrable martingale with respect to $\left(\mathcal{G}_{t}^{A_{m}}\right)$ and its predictable quadratic variation is given by $\langle G(m)\rangle=\int_{0}^{t} \widehat{\sigma}_{s} \widehat{\sigma}_{s}^{*} \mathrm{~d} s+\left(\widetilde{\delta}^{2} 1_{A_{m}^{c} \cap\{\widetilde{\delta} \mid \leq 1\}}\right) \star v$, and $G(m)_{S_{r}-}=G(m)_{S_{r}}$. Since $S_{r}$ is $\mathcal{G}_{0}^{A_{m}}$-measurable, $\left(S_{r}-h\right)_{+}$is a $\left(\mathcal{G}_{t}^{A_{m}}\right)$-stopping time for every $h \geq 0$. Therefore, the Lenglart inequality implies that

$$
\begin{aligned}
& P\left(\sup _{0 \leq h \leq h_{0}}\left|h_{0}^{-1 / 2}\left\{G(m)_{S_{r}}-G(m)_{\left(S_{r}-h\right)_{+}}\right\}\right|^{2}>K\right) \\
& \quad \leq \frac{K^{\prime}}{K}+P\left(h_{0}^{-1}\left|\langle G(m)\rangle_{S_{r}}-\langle G(m)\rangle_{\left(S_{r}-h_{0}\right)_{+}}\right|>K^{\prime}\right)
\end{aligned}
$$

for any $K, K^{\prime}>0$, and thus a standard localization argument yields $\sup _{0 \leq h \leq h_{0}} \mid G(m)_{S_{r}}-$ $G(m)_{\left(S_{r}-h\right)_{+}} \mid=\mathrm{O}_{p}\left(\sqrt{h_{0}}\right)$. This completes the proof of the lemma.

Lemma 6.4. Under [SA1], $N_{S_{r}-}^{n}-N_{S_{r}^{\dagger}}^{n}=k_{n}+\mathrm{o}_{p}\left(n^{(1 / 2)-\alpha^{\prime}}\right)$ for any $\alpha^{\prime} \in\left(0,\left(\xi-\kappa-\frac{1}{2}\right) \wedge(\xi-\right.$ $\left.\left.\frac{3}{4}\right) \wedge \varpi\right)$.

Proof. Since $\alpha^{\prime}<\xi-\kappa-\frac{1}{2}$, by [SA1], we have

$$
N_{S_{r}-}^{n}-N_{S_{r}^{\dagger}}^{n}=\sum_{p=1}^{N_{S_{r}-1}^{n}+1} \frac{E\left[n\left|I_{p}\right| \mid \mathcal{G}_{T_{p-1}}^{(0)}\right]}{G_{T_{p-1}}^{n}} 1_{\left\{T_{p-1}>S_{r}^{\dagger}\right\}}+\mathrm{o}_{p}\left(n^{1 / 2-\alpha^{\prime}}\right) .
$$

In particular, from this expression and [SA1], we deduce $N_{S_{r}-}^{n}-N_{S_{r}^{\dagger}}^{n}=\mathrm{O}_{p}(\sqrt{n} \log n)$. Therefore, [SA1] yields

$$
\sum_{p=1}^{N_{S_{r}-}^{n}+1} E\left[\left.\left|n^{\alpha^{\prime}-1 / 2} \frac{n\left|I_{p}\right|}{G_{T_{p-1}}^{n}} 1_{\left\{T^{p-1}>S_{r}^{\dagger}\right\}}\right|^{2} \right\rvert\, \mathcal{G}_{T_{p-1}}^{(0)}\right]=n^{1+2 \alpha^{\prime} \bar{r}_{n}^{2}} \sum_{p=1}^{N_{S r_{r}-}^{n}+1} \frac{1}{G_{T_{p-1}}^{n}} 1_{\left\{T^{p-1}>S_{r}^{\dagger}\right\}}=\mathrm{o}_{p}(1),
$$

hence Lemma 2.3 of [18] implies that

$$
N_{S_{r}-}^{n}-N_{S_{r}^{\dagger}}^{n}=n \sum_{p=1}^{N_{S_{r}-}^{n}+1} \frac{\left|I_{p}\right|}{G_{T_{p-1}}^{n}} 1_{\left\{T^{p-1}>S_{r}^{\dagger}\right\}}+\mathrm{o}_{p}\left(n^{(1 / 2)-\alpha^{\prime}}\right) .
$$

Now [SA1], Lemma 6.3 and the fact that $\alpha^{\prime}<\varpi \wedge \frac{1}{4}$ yield

$$
\begin{aligned}
n \sum_{p=1}^{N_{S_{r}-1}^{n}+1}\left(\frac{\left|I_{p}\right|}{G_{T_{p-1}}^{n}}-\frac{\left|I_{p}\right|}{G_{\underline{S}_{r}}^{n}}\right) 1_{\left\{T^{p-1}>S_{r}^{\dagger}\right\}} & =n \sum_{p=1}^{N_{S_{r}-}^{n}+1}\left(\frac{\left|I_{p}\right|}{G_{T_{p-1}}}-\frac{\left|I_{p}\right|}{G_{\underline{S}_{r}}}\right) 1_{\left\{T^{p-1}>S_{r}^{\dagger}\right\}}+\mathrm{o}_{p}\left(n^{(1 / 2)-\alpha^{\prime}}\right) \\
& =\mathrm{o}_{p}\left(n^{(1 / 2)-\alpha^{\prime}}\right)
\end{aligned}
$$

thus we have

$$
N_{S_{r}-}^{n}-N_{S_{r}^{\dagger}}^{n}=n \sum_{p=1}^{N_{S_{r}-}^{n}+1} \frac{\left|I_{p}\right|}{G_{\underline{S}_{r}}^{n}} 1_{\left\{T^{p-1}>S_{r}^{\dagger}\right\}}+\mathrm{o}_{p}\left(n^{(1 / 2)-\alpha^{\prime}}\right)=k_{n} \frac{G_{\underline{S}_{r}}^{n} \wedge \log n}{G_{\underline{S}_{r}}^{n}}+\mathrm{o}_{p}\left(n^{(1 / 2)-\alpha^{\prime}}\right) .
$$

Since $\lim _{n} P\left(G_{\underline{S}_{r}}^{n}>\log n\right)=0$, we obtain the desired result.

Now we proceed to the main body of the proof of Proposition 6.6. Denote by $\Omega_{n}(m)$ the set on which $\left|S_{r_{1}}-S_{r_{2}}\right|>\left(k_{n} / n\right) \log n$ for any $r_{1}, r_{2} \in \mathcal{R}_{m}$ such that $r_{1} \neq r_{2}$ and $S_{r_{1}}, S_{r_{2}}<\infty$. Since $S_{r_{1}} \neq S_{r_{2}}$ if $r_{1} \neq r_{2}$ and $S_{r_{1}}, S_{r_{2}}<\infty$, we have $P\left(\Omega_{n}(m)\right) \rightarrow 1$ as $n \rightarrow \infty$.

Lemma 6.5. Under $[\mathrm{SA} 2], E\left[\sup _{\underline{S}_{r} \leq s<S_{r}}\left\|\sigma_{s}-\sigma_{\underline{S}_{r}}\right\|^{2} ; \Omega_{n}(m)\right] \lesssim\left(k_{n} / n\right) \log n$.

Proof. Since no jump of the Poisson process $1_{A_{m}} \star \mu$ occurs in $\left[\underline{S}_{r}, S_{r}\right)$ on the set $\Omega_{n}(m)$, we have $\sigma_{s}=\sigma(m)_{s}$ for every $s \in\left[\underline{S}_{r}, S_{r}\right)$ on this set, where

$$
\left\{\begin{array}{l}
\sigma(m)_{s}=\sigma_{0}+\int_{0}^{s} \widetilde{b}(m)_{u} \mathrm{~d} u+\int_{0}^{s} \widetilde{\sigma}_{u} \mathrm{~d} W_{u}+\left(\widetilde{\delta} 1_{A_{m}^{c}}\right) \star(\mu-v)_{s} \\
\widetilde{b}(m)_{u}=\widetilde{b}_{u}-\int_{A_{m} \cap\{z:|\widetilde{\delta}(u, z)| \leq 1\}} \widetilde{\delta}(u, z) \lambda(\mathrm{d} z) .
\end{array}\right.
$$

On the other hand, by Proposition 2.1.10 of [31] we have $E\left[\sup _{\underline{S}_{r} \leq s<S_{r}}\left\|\sigma(m)_{s}-\sigma(m) \underline{\underline{S}}_{r}\right\|^{2}\right] \lesssim$ $\frac{k_{n}}{n} \log n$, which implies the desired result.

Proof of Proposition 6.6. Throughout the proof, we fix a constant $\alpha^{\prime}$ such that $1-\xi<\alpha^{\prime}<$ $\left(\xi-\kappa-\frac{1}{2}\right) \wedge\left(\xi-\frac{3}{4}\right) \wedge \varpi$. Such an $\alpha^{\prime}$ exists due to (6.1).

First we prove the first equation of (6.7). Set $\Omega_{n}=\left\{N_{S_{r-}}^{n}-k_{n}+1 \geq 0\right\}$. By the Lipschitz continuity of $g$,

$$
\begin{aligned}
\eta_{-}(n, r)^{k} & =\frac{n^{1 / 4}}{\psi_{2}} \sum_{p=N_{S_{r}-}^{n}-k_{n}+2}^{N_{S_{r}-}^{n}}\left(\phi_{g, g}\right)_{N_{S_{r}-}^{n}+1-p}^{n} M^{k}\left(I_{p}\right)+\mathrm{o}_{p}(1) \\
& =\frac{n^{1 / 4}}{\psi_{2}} \sum_{w=1}^{k_{n}-1}\left(\widetilde{\phi_{g, g}}\right)_{w}^{n} M^{k}\left(I_{N_{S_{r}-}+1-k_{n}+w}\right)+\mathrm{o}_{p}(1)
\end{aligned}
$$

on $\Omega_{n}$. On the other hand, noting that we have

$$
\sup \left\{\left\|M_{s}-M_{r}\right\|:|s-r| \leq h, s, r \in[0, t]\right\}=\mathrm{O}_{p}(\sqrt{h}|\log h|)
$$

as $h \downarrow 0$ for any $t>0$ due to a representation of a continuous local martingale with Brownian motion and Lévy's theorem on the uniform modulus of continuity of Brownian motion, summation by parts, (6.2) and Lemma 6.4 imply that

$$
\begin{aligned}
& n^{1 / 4} \sum_{w=1}^{k_{n}-1}\left(\widetilde{\phi_{g, g}}\right)_{w}^{n}\left\{M^{k}\left(I_{N_{S_{r}-}}+1-k_{n}+w\right)-M^{k}\left(I_{i\left(S_{r}^{\dagger}\right)^{n}+w}\right)\right\} \\
& =n^{1 / 4} \sum_{w=1}^{k_{n}-2}\left\{\left(\widetilde{\phi_{g, g}}\right)_{w}^{n}-\left(\widetilde{\phi_{g, g}}\right)_{w+1}^{n}\right\}\left(M_{T_{N_{S_{r}-}+1-k_{n}+w}^{k}}^{k}-M_{T_{i\left(S_{r}^{\dagger}\right)^{n}+w}^{k}}^{k}\right) \\
& \quad+n^{1 / 4}\left(\widetilde{\phi_{g, g}}\right)_{k_{n}-1}^{n}\left(M_{T_{N_{S_{r}-}^{n}}^{n}}^{k}-M_{T_{i\left(S_{r}\right)^{n}+k_{n}-1}^{k}}^{k}\right)-n^{1 / 4}\left(\widetilde{\phi_{g, g}}\right)_{1}^{n}\left(M_{T_{S_{S_{r}-}^{n}+1-k_{n}}^{k}}^{k}-M_{T_{i\left(S_{r}\right) n}^{n} n}^{k}\right) \\
& =\mathrm{o}_{p}\left(n^{1 / 4} \sqrt{n^{1 / 2-\alpha^{\prime}-\xi} \log n}\right)=\mathrm{o}_{p}\left(n^{\left(1-\xi-\alpha^{\prime}\right) / 2} \sqrt{\log n}\right)=\mathrm{o}_{p}(1)
\end{aligned}
$$

on $\Omega_{n}$. Since $\lim _{n} P\left(\Omega_{n}\right)=1$, we conclude that

$$
\eta_{-}(n, r)^{k}=\frac{n^{1 / 4}}{\psi_{2}} \sum_{w=1}^{k_{n}-1}\left(\widetilde{g_{g, g}}\right)_{w}^{n} M^{k}\left(I_{i\left(S_{r}^{\dagger}\right)+w}\right)+\mathrm{o}_{p}(1) .
$$

Next, noting that $W$ is a $d^{\prime}$-dimensional $\left(\mathcal{G}_{t}^{(0)}\right)$-Brownian motion (recall that $\left(\mathcal{G}_{t}^{(0)}\right)$ is the smallest filtration containing $\left(\mathcal{F}_{t}^{(0)}\right)$ such that $\mathcal{G}_{0}^{(0)}$ contains the $\sigma$-field generated by $\mu$ ), we have

$$
\begin{aligned}
& n^{1 / 4} \sum_{w=1}^{k_{n}-1}\left(\widetilde{\phi_{g, g}}\right)_{w}^{n} M^{k}\left(I_{i\left(S_{r}^{\dagger}\right)+w}\right)-\sum_{j=1}^{d^{\prime}} \sigma_{\underline{S_{r}}}^{k j} z_{r-}^{n, j} \\
& \quad=n^{1 / 4} \sum_{j=1}^{d^{\prime}} \sum_{w=1}^{k_{n}-1}\left(\widetilde{\phi_{g, g}}\right)_{w}^{n} \int_{T_{i\left(S_{r}^{\dagger}\right)^{n}+w-1}^{T}}^{T_{i\left(S_{r}^{\dagger}\right)^{n}+w}}\left(\sigma_{s}^{k j}-\sigma_{\underline{S}_{r}}^{k j}\right) \mathrm{d} W_{s}^{j},
\end{aligned}
$$

hence the Lenglart inequality implies that it is enough to show that

$$
\Delta_{n}:=\sqrt{n} \sum_{j=1}^{d^{\prime}} \sum_{w=1}^{k_{n}-1}\left|\left(\widetilde{\phi_{g}, g}\right)_{w}^{n}\right|^{2} E\left[\int_{T_{i\left(S_{r}^{\dagger}\right)^{n}+w-1}}^{T_{i\left(S_{r}^{\dagger}\right)^{n}+w}\left(\sigma_{s}^{k j}-\sigma_{\underline{S}_{r}}^{k j}\right)^{2} \mathrm{~d} s \mid \mathcal{G}_{T_{i\left(S_{r}^{\dagger}\right)^{n}+w-1}^{(0)}}^{()^{p}} 0 . . \text {. } \rightarrow^{p} 0 .}\right.
$$

Set

$$
\Delta_{n}^{\prime}=\sqrt{n} \sum_{j=1}^{d^{\prime}} \sum_{w=1}^{k_{n}-1}\left|\left(\widetilde{\phi_{g, g}}\right)_{w}^{n}\right|^{2} E\left[\int_{T_{i\left(S_{r}^{\dagger}\right)^{n}+w-1}}^{T_{i\left(S S_{i}^{\dagger}\right)^{n}+w}\left(\sigma_{s}^{k j}-\sigma_{\underline{S}_{r}}^{k j}\right)^{2} \mathrm{~d} s \mid \mathcal{G}_{T_{i\left(S_{r}^{\dagger}\right)^{n}+w-1}^{(0)}}^{(1)} 1_{\left\{T_{i\left(S_{r}^{\dagger}\right)^{\dagger}+w-1} \leq S_{r}\right\}} . . . . ~}\right.
$$

Then, since $\Omega_{n}(m) \in \mathcal{G}_{0}^{(0)}$, it holds that

$$
E\left[\Delta_{n}^{\prime} ; \Omega_{n}(m)\right] \lesssim \sqrt{n} \sum_{j=1}^{d^{\prime}} E\left[\int_{T_{i\left(S_{r}^{\dagger}\right)}^{n}}^{T_{i\left(S_{r}\right)^{n}}}\left(\sigma_{s}^{k j}-\sigma_{S_{r}^{\dagger}}^{k j}\right)^{2} \mathrm{~d} s ; \Omega_{n}(m)\right] .
$$

Now, Lemma 6.5, the boundedness of $\sigma$ and (6.2) imply that

$$
E\left[\int_{\left.T_{i\left(S_{r}^{\dagger}\right)^{n}}^{T_{i\left(S_{r}\right)^{n}}}\left(\sigma_{s}^{k j}-\sigma_{S_{r}^{\dagger}}^{k j}\right)^{2} \mathrm{~d} s ; \Omega_{n}(m)\right] \lesssim \frac{k_{n}}{n}(\log n) E\left[\sup _{\underline{S_{r}} \leq s<S_{r}}\left(\sigma_{s}^{k j}-\sigma_{S_{r}^{\dagger}}^{k j}\right)^{2} ; \Omega_{n}(m)\right]+\bar{r}_{n} \lesssim \bar{r}_{n}, ~, ~ . ~}^{\text {, }}\right.
$$

hence we obtain $E\left[\Delta_{n}^{\prime} ; \Omega_{n}(m)\right] \lesssim \sqrt{n} \bar{r}_{n}=\mathrm{o}(1)$. Therefore, the equation $\lim _{n} P\left(\Omega_{n}(m)\right)=1$ and the Chebyshev inequality yield $\Delta_{n}^{\prime}=\mathrm{o}_{p}(1)$. On the other hand, the boundedness of $\sigma$, (6.2) and Lemma 6.4 imply that $\left|\Delta_{n}-\Delta_{n}^{\prime}\right|=\mathrm{o}_{p}\left(n^{1-\xi-\alpha^{\prime}}\right)=\mathrm{o}_{p}(1)$. Consequently, we obtain $\Delta_{n}=\mathrm{o}_{p}(1)$ and the first equation of (6.7) has been proved. On the other hand, noting that $N_{S_{r}}^{n}-N_{S_{r}-}^{n} \leq 1$ and $S_{r}$ is an $\left(\mathcal{F}_{t}^{(0)}\right)$-stopping time, the second equation of (6.7) can be shown in a similar (and simpler) manner.

Next, we prove the first equation of (6.8). By the (piecewise) Lipschitz continuity of $g$ and $g^{\prime}$, we have on $\Omega_{n}$

$$
\eta_{-}^{\prime}(n, r)^{k}=-\frac{n^{1 / 4}}{\psi_{2} k_{n}} \sum_{p=\left(N_{S_{r}-}^{n}-k_{n}+2\right)_{+}}^{N_{S_{r-}}^{n}}\left(\phi_{g^{\prime}, g}\right)_{N_{S_{r}-}^{n}+1-p}^{n} \epsilon_{\tau_{p}^{k}}^{k}+\mathrm{o}_{p}(1)
$$

Moreover, by Lemma 6.4, [SA3] and the Lipschitz continuity of $\phi_{g^{\prime}, g}$ we have

$$
\begin{aligned}
& E_{0}\left[\left|\frac{n^{1 / 4}}{k_{n}}\left\{\sum_{p=\left(N_{S_{r}-}^{n}-k_{n}+2\right)_{+}}^{N_{S_{r}-}^{n}}\left(\phi_{g^{\prime}, g}\right)_{N_{S_{r}-}^{n}+1-p}^{n} \epsilon_{\tau_{p}^{k}}^{k}-\sum_{p=i\left(S_{r}^{\dagger}\right)^{n}+1}^{i\left(S_{r}^{\dagger}\right)^{n}+k_{n}-1}\left(\phi_{g^{\prime}, g}\right)_{i\left(S_{r}^{\dagger}\right)^{n}+k_{n}-p}^{n} \epsilon_{\tau_{p}^{k}}^{k}\right)\right|^{2}\right] \\
& \quad=\mathrm{O}_{p}\left(n^{-\alpha^{\prime}}\right)
\end{aligned}
$$

on $\Omega_{n}$. Since $\lim _{n} P\left(\Omega_{n}\right)=1$, we conclude that

$$
\eta_{-}^{\prime}(n, r)^{k}=-\frac{n^{1 / 4}}{\psi_{2} k_{n}} \sum_{p=i\left(S_{r}^{\dagger}\right)^{n}+1}^{i\left(S_{r}^{\dagger}\right)^{n}+k_{n}-1}\left(\phi_{g^{\prime}, g}\right)_{i\left(S_{r}^{\dagger}\right)^{n}+k_{n}-p}^{n} \epsilon_{\tau_{p}^{k}}^{k}+\mathrm{o}_{p}(1)=z_{r-}^{\prime n, k}+\mathrm{o}_{p}(1) .
$$

Similarly, we can prove the second equation of (6.8).

### 6.2.2. An auxiliary local stable convergence result

In this subsection, we prove an auxiliary local stable convergence result corresponding to Lemma 16.3.7 of [31]. The proof is close to that of the aforementioned lemma, but there is a difference due to the additional randomness coming from the sampling times. Furthermore, we can also simplify some parts of the proof because it is sufficient for our purpose to prove a simpler consequence than that of the aforementioned lemma. For these reasons, we give a complete proof.

The following lemma is a direct consequence of the Skorokhod representation theorem, so we omit the proof:

Lemma 6.6. Let $\left(f_{n}\right)$ be a sequence of real-valued functions on $\mathbb{R}^{D}$ such that there exists a constant $K$ satisfying $\left|f_{n}(x)\right| \leq K$ and $\left|f_{n}(x)-f_{n}(y)\right| \leq K\|x-y\|$ for all $x, y \in \mathbb{R}^{D}$ and every $n$. If a sequence $\left(x_{n}\right)$ of $\mathbb{R}^{D}$-valued random variables converges in law to a variable $x$, then $E\left[f_{n}\left(x_{n}\right)\right]-E\left[f_{n}(x)\right] \rightarrow 0$.

The following lemma is the main result of this subsection. We denote by $\mathfrak{N}_{D}$ the $D$ dimensional standard normal distribution.

Lemma 6.7. Assume that [SA1], [SA3] and [A4](ii) are satisfied. Suppose that for each $n$ there is a $\left(\mathcal{G}_{t}^{(0)}\right)$-stopping time $\rho_{n}$. Suppose also that there is a finite-valued variable $\rho$ such
that $\rho_{n} \rightarrow \rho$ as $n \rightarrow \infty$ and one of the following two condition is satisfied:
(1) $\rho>0, P\left(T_{i\left(\rho_{n}\right)^{n}+k_{n}-\left\lfloor n^{\beta}\right\rfloor}<\rho\right) \rightarrow 1$ as $n \rightarrow \infty$ for some $\beta \in(0, \xi-1 / 2)$,
$\left.\begin{array}{l}\text { in which case we set } G_{(\rho)}=G_{\rho-}, \widetilde{v}_{(\rho)}=\widetilde{v}_{\rho-} \text { and } \mathcal{G}_{(\rho)}^{(0)}=\mathcal{G}_{\rho-}^{(0)} \text {, } \\ \text { (2) } \quad \rho_{n} \geq \rho \text { for all } n \text {, in which case we set } G_{(\rho)}^{(0)}=G_{\rho}, \widetilde{v}_{(\rho)}=\widetilde{v}_{\rho} \text { and } \mathcal{G}_{(\rho)}^{(0)}={ }_{\rho}\end{array}\right\}$.
Let $\phi_{1}$ and $\phi_{2}$ be continuous real-valued functions $\phi_{1}$ and $\phi_{2}$ on $[0,1]$. Then, for any $\mathcal{F}$ measurable bounded variable $U$ and any bounded Lipschitz function $f$ on $\mathbb{R}^{d^{\prime}+d}$ we have

$$
\begin{align*}
& E\left[U f\left(L_{n}, L_{n}^{\prime}\right) \mid \mathcal{G}_{\rho_{n}}\right]  \tag{6.10}\\
& \quad \rightarrow^{p} E\left[U \int f\left(\left\|\phi_{1}\right\| \sqrt{\theta G_{(\rho)}} x,\left\|\phi_{2}\right\| \sqrt{\theta^{-1}} \widetilde{v}_{(\rho)} y\right) \mathfrak{N}_{d^{\prime}}(\mathrm{d} x) \mathfrak{N}_{d}(\mathrm{~d} y) \mid \mathcal{G}_{(\rho)}^{(0)}\right],
\end{align*}
$$

where $L_{n}=L\left(\phi_{1}, \rho_{n}\right)_{n}, L_{n}^{\prime}=L^{\prime}\left(\phi_{2}, \rho_{n}\right)_{n}$ and $\left\|\phi_{j}\right\|^{2}=\int_{0}^{1} \phi_{j}(x)^{2} \mathrm{~d} x$ for $j=1,2$.
Proof. Step 1. For $k, l=1, \ldots, d$ we set $D_{n}^{k l}=\frac{1}{k_{n}} \sum_{w=1}^{k_{n}-1}\left|\left(\phi_{2}\right)_{w}^{n}\right|^{2} 1_{\left\{\tau_{i\left(\rho_{n}\right)^{n}+w}^{k}=\tau_{i\left(\rho_{n}\right)^{n}+w}^{l}\right\}}$. We begin by proving $D_{n}^{k l} \rightarrow^{p}\left\|\phi_{2}\right\|^{2} \chi_{(\rho)}^{k l}$, where we set $\chi_{(\rho)}^{k l}=\chi_{\rho-}^{k l}$ in case (1) and $\chi_{(\rho)}^{k l}=\chi_{\rho}^{k l}$ in case (2). Since $i\left(\rho_{n}\right)$ is a $\left(\mathcal{G}_{T_{p}}^{(0)}\right)_{p=0}^{\infty}$-stopping time, [SA1] and Lemma 2.3 of [18] yield $D_{n}^{k l}=\frac{1}{k_{n}} \sum_{w=1}^{k_{n}-1}\left|\left(\phi_{2}\right)_{w}^{n}\right|^{2} \chi_{T_{i\left(\rho_{n}\right)^{n}+w-1}^{k l}}+\mathrm{o}_{p}(1)$. Since $\chi^{k l}$ is cadlag, (6.9) implies that $D_{n}^{k l}=$ $\left\|\phi_{2}\right\|^{2} \chi_{(\rho)}^{k l}+\mathrm{o}_{p}(1)$.

Step 2. From step 1, by considering an appropriate subsequence if necessary, without loss of generality we may assume that there is a subset $\Omega_{0}$ of $\Omega^{(0)}$ such that $P^{(0)}\left(\Omega_{0}\right)=1$ and $D_{n}^{k l}\left(\omega^{(0)}\right) \rightarrow\left\|\phi_{2}\right\|^{2} \chi_{(\rho)}^{k l}\left(\omega^{(0)}\right)$ for all $\omega^{(0)} \in \Omega_{0}$.

Step 3. Fix $\omega^{(0)} \in \Omega_{0}$, and consider the probability space $\left(\Omega^{(1)}, \mathcal{F}^{(1)}, Q_{0}\right)$, where $Q_{0}(\cdot)=$ $Q\left(\omega^{(0)}, \cdot\right)$. Our aim in this step is to show that under $Q_{0}$

$$
\begin{equation*}
L_{n}^{\prime} \rightarrow^{d}\left\|\phi_{2}\right\| \sqrt{\theta^{-1}} \widetilde{v}_{(\rho)}\left(\omega^{(0)}\right) \zeta^{\prime} \tag{6.11}
\end{equation*}
$$

where $\zeta^{\prime}$ is a standard $d$-dimensional normal variable independent of $\mathcal{F}$.
For each $w=1, \ldots, k_{n}-1$ we define the $d$-dimensional variable $y_{w}^{n}=\left(y_{w}^{n, k}\right)_{1 \leq k \leq d}$ by

$$
y_{w}^{n, k}=\frac{n^{1 / 4}}{k_{n}}\left(\phi_{2}\right)_{w}^{n} \epsilon_{\tau_{i\left(\rho_{n}\right)^{n}\left(\omega^{(0)}\right)+w}^{k}}^{k} .
$$

Then $y_{1}^{n}, \ldots, y_{k_{n}-1}^{n}$ are independent under $Q_{0}$ and we have $L_{n}^{\prime}\left(\omega^{(0)}, \cdot\right)=\sum_{w=1}^{k_{n}-1} y_{w}^{n}$. Moreover, by [SA3] we have

$$
\begin{aligned}
E_{Q_{0}}\left(y_{w}^{n}\right) & =0, \quad E_{Q_{0}}\left(\left\|y_{w}^{n}\right\|^{4}\right) \lesssim k_{n}^{-2}, \quad \sum_{w=1}^{k_{n}-1} E_{Q_{0}}\left(\left\|y_{w}^{n}\right\|^{4}\right) \rightarrow 0, \\
\sum_{w=1}^{k_{n}-1} E_{Q_{0}}\left(y_{w}^{n, k} y_{w}^{n, l}\right) & =\frac{n^{1 / 2}}{k_{n}^{2}} \sum_{w=1}^{k_{n}-1}\left|\left(\phi_{2}\right)_{w}^{n}\right|^{2} \Upsilon^{k l}\left(\omega^{(0)}\right)_{\tau_{i\left(\rho_{n}\right)^{n}\left(\omega^{(0)}\right)+w}^{k}} 1_{\left\{\tau_{i\left(\rho_{n}\right)^{n}\left(\omega^{(0)}\right)+w}^{k}=\tau_{i\left(\rho_{n}\right)^{n}\left(\omega^{(0)}\right)+w}^{l}\right\}} \cdot
\end{aligned}
$$

Since $\Upsilon$ is cadlag, (6.9) and the fact that $n^{1 / 2} / k_{n} \rightarrow \theta^{-1}$ yield $\sum_{w=1}^{k_{n}-1} E_{Q_{0}}\left(y_{w}^{n, k} y_{w}^{n, l}\right)=$ $\theta^{-1} \Upsilon_{(\rho)}^{k l}\left(\omega^{(0)}\right) D_{n}^{k l}\left(\omega^{(0)}\right)+o_{p}(1)$, where we set $\Upsilon_{(\rho)}^{k l}=\Upsilon_{\rho-}^{k l}$ in case (1) and $\Upsilon_{(\rho)}^{k l}=\Upsilon_{\rho}^{k l}$ in case (2). Since $\omega^{(0)} \in \Omega^{(0)}$, this implies that

$$
\sum_{w=1}^{k_{n}-1} E_{Q_{0}}\left(y_{w}^{n, k} y_{w}^{n, l}\right) \rightarrow^{p}\left\|\phi_{2}\right\|^{2} \theta^{-1} \Upsilon_{(\rho)}^{k l}\left(\omega^{(0)}\right) \chi_{(\rho)}^{k l}\left(\omega^{(0)}\right)
$$

Now a standard central limit theorem on row-wise independent triangular arrays of infinitesimal variables (e.g., Theorem 2.2.14 of [31]) yields (6.11).

Step 4. In this step, we shall show the following convergence for $L_{n}$ :

$$
\begin{equation*}
L_{n} \rightarrow^{d} \sqrt{\Phi_{22} \theta G_{(\rho)}} \zeta \tag{6.12}
\end{equation*}
$$

where $\zeta$ is a standard $d^{\prime}$-dimensional normal variable independent of $\mathcal{F}$. Unlike step 3, here the limiting variable is mixed normal, so we cannot rely on the standard central limit theorem used in step 3. Instead, we use the classic mixed normal limit theorem of [20].

Fix $u \in \mathbb{R}^{d^{\prime}}$ arbitrarily and set $y(u)_{w}^{n}=n^{1 / 4}\left(\phi_{1}\right)_{w}^{n} u^{*} W\left(I_{i\left(\rho_{n}\right)^{n}+w}\right)$ for each $w=1, \ldots, k_{n}-1$.
 $G_{-}$do not vanish, it suffices to verify the following four conditions according to [20] and the Cramér-Wold method:

$$
\begin{align*}
E\left[\max _{1 \leq w \leq k_{n}-1}\left|y(u)_{w}^{n}\right|^{2}\right] & \rightarrow 0  \tag{6.13}\\
\sum_{w=1}^{k_{n}-1}\left|y(u)_{w}^{n}\right|^{2}-\|u\|^{2}\left\|\phi_{1}\right\|^{2} \theta G_{\rho_{n}} & \rightarrow^{p} 0,  \tag{6.14}\\
\|u\|^{2}\left\|\phi_{1}\right\|^{2} \theta G_{\rho_{n}} & \rightarrow^{p}\|u\|^{2}\left\|\phi_{1}\right\|^{2} \theta G_{(\rho)},  \tag{6.15}\\
\sum_{w=1}^{k_{n}-1} \mid E\left[y(u)_{w}^{n} \mid \mathcal{G}_{\left.T_{i\left(\rho_{n}\right)^{n}+w-1}\right]}\right] & \rightarrow^{p} 0 . \tag{6.16}
\end{align*}
$$

Equation (6.13) follows from (6.2) and Lévy's theorem on the uniform modulus of continuity of Brownian motion. Next, [SA1] and Lemma 2.3 of [18] imply that

$$
\sum_{w=1}^{k_{n}-1}\left|y(u)_{w}^{n}\right|^{2}=\|u\|^{2} \sqrt{n} \sum_{w=1}^{k_{n}-1}\left(\phi_{1}\right)_{w}^{n}\left|I_{i\left(\rho_{n}\right)^{n}+w}\right|+\mathrm{o}_{p}(1)=\frac{\|u\|^{2}}{\sqrt{n}} \sum_{w=1}^{k_{n}-1}\left(\phi_{1}\right)_{w}^{n} G_{T_{i\left(\rho_{n}\right)^{n}+w-1}}+\mathrm{o}_{p}(1)
$$

hence we obtain (6.14) because $G$ is cadlag. Finally, the fact that $G$ is cadlag and (6.9) yield (6.15), while we have $E\left[y(u)_{w}^{n} \mid \mathcal{G}_{T_{i\left(\rho_{n}\right)^{n}+w-1}}\right]=0$ because $W$ is a $d^{\prime}$-dimensional $\left(\mathcal{F}_{t}\right)$ Brownian motion independent of $\mathcal{G}$, hence (6.16) holds true.

Step 5. We denote by $\Psi_{n}(U)$ and $\Psi(U)$ the left-hand and right-hand sides of (6.10), respectively. In this step, we show that it is enough to prove

$$
\begin{equation*}
\Psi_{n}(1) \rightarrow^{p} \Psi(1) \tag{6.17}
\end{equation*}
$$

In fact, assume this, and take an arbitrary bounded variable $U$. We consider the cadlag version of the bounded martingale $U_{t}=E\left(U \mid \mathcal{G}_{t}^{(0)}\right)$.

First, suppose that we are in case (1). Set $\underline{k}_{n}=k_{n}-\left\lfloor n^{\beta}\right\rfloor$ and define the $d^{\prime}$-dimensional variable $\underline{L}_{n}=\left(\underline{L}_{n}^{j}\right)_{1 \leq j \leq d^{\prime}}$ and the $d$-dimensional variable $\underline{L}_{n}^{\prime}=\left(\underline{L}_{n}^{\prime}\right)_{1 \leq k \leq d}$ by

$$
\underline{L}_{n}^{j}=n^{1 / 4} \sum_{w=1}^{\underline{k}_{n}}\left(\phi_{1}\right)_{w}^{n} W^{j}\left(I_{i(\rho)^{n}+w}\right), \quad \underline{L}_{n}^{\prime k}=\frac{n^{1 / 4}}{k_{n}} \sum_{w=1}^{\underline{k}_{n}}\left(\phi_{2}\right)_{w}^{n} \epsilon_{\tau_{i(\rho)^{n}+w}^{k}}^{k}
$$

Then, since $E\left[\left\|L_{n}-\underline{L}_{n}\right\|^{2}\right] \lesssim \sqrt{n} n^{\beta} \bar{r}_{n}$ and $E\left[\left\|L_{n}^{\prime}-\underline{L}_{n}^{\prime}\right\|^{2}\right] \lesssim \sqrt{n} k_{n}^{-2} n^{\beta}$, by the boundedness of $U$ and the Lipschitz continuity of $f$ it holds that $\Psi_{n}(U)-\underline{\Psi}_{n}(U) \rightarrow^{p} 0$ and $\Psi_{n}(1)-$ $\underline{\Psi}_{n}(1) \rightarrow^{p} 0$, where $\underline{\Psi}_{n}(U)=E\left[U f\left(\underline{L}_{n}, \underline{L}_{n}^{\prime}\right) \mid \mathcal{G}_{\rho_{n}}\right]$. In particular, to prove (6.17) it is enough to show that $\underline{\Psi}_{n}(U) \rightarrow^{p} \Psi(U)$. Now, since both $G_{\rho-}$ and $\widetilde{v}_{\rho-}$ are $\mathcal{G}_{\rho-}^{(0)}$-measurable, we have $\Psi(U)=U_{\rho-} \Psi(1)$ because $U_{\rho-}=E\left[U \mid \mathcal{G}_{\rho-}^{(0)}\right]$. Also, $f\left(\underline{L}_{n}, \underline{L}_{n}^{\prime}\right)$ in restriction to the set $\Omega_{n}=$ $\left\{\rho>T_{i\left(\rho_{n}\right)^{n}+\underline{k}_{n}}\right\}$ is $\mathcal{G}_{\rho--}$-measurable, so $\underline{\Psi}_{n}(U)=\underline{\Psi}_{n}\left(U_{\rho-}\right)$ on $\Omega_{n}$. We also obviously have $\underline{\Psi}_{n}\left(U_{\rho_{n}}\right)=\underline{\Psi}_{n}(1) U_{\rho_{n}} \rightarrow^{p} \Psi(1) U_{\rho-}$ by (6.17), $\Psi_{n}(1)-\underline{\Psi}_{n}(1) \rightarrow^{p} 0$ and $U_{\rho_{n}} \rightarrow U_{\rho-}$, while $P\left(\Omega_{n}\right) \rightarrow 1$ by assumption. Now, since $E\left[\left|\underline{\Psi}_{n}\left(U_{\rho_{n}}\right)-\underline{\Psi}_{n}\left(U_{\rho-}\right)\right|\right] \leq\|f\|_{\infty} E\left[\left|U_{\rho_{n}}-U_{\rho-}\right|\right] \rightarrow 0$ by the boundedness of $f$ and $U, U_{\rho_{n}} \rightarrow U_{\rho-}$ on $\Omega_{n}$ and the fact that $P\left(\Omega_{n}\right) \rightarrow 1$, we obtain the desired result.

Next, suppose that we are in case (2). Then $\Psi(U)=U_{\rho} \Psi(1)$ because $\Psi(1)$ is $\mathcal{G}_{\rho}^{(0)}$-measurable, and also $\Psi_{n}\left(U_{\rho}\right)=U_{\rho} \Psi_{n}(1)$ because $\rho_{n} \geq \rho$. Moreover, setting $\rho_{n}^{\prime}=\rho_{n}+k_{n} \bar{r}_{n}, \Psi_{n}(1)$ is $\mathcal{G}_{\rho_{n}^{\prime}}-$ measurable due to (6.2), so $\Psi_{n}(U)=U_{\rho_{n}^{\prime}} \Psi_{n}(1)$. Since $\rho_{n}^{\prime} \rightarrow \rho$ and $\rho_{n}^{\prime}>\rho$, we have $U_{\rho_{n}^{\prime}} \rightarrow U_{\rho}$, and the same arguments as above shows that $\Psi_{n}\left(U_{\rho_{n}^{\prime}}\right)-\Psi_{n}\left(U_{\rho}\right) \rightarrow^{p} 0$, thus the desired result is obtained.

Step 6. Now we finish the proof by proving the convergence (6.17). First, for each $\omega^{(0)} \in \Omega_{0}$ define the function $h_{\omega^{(0)}}^{n}$ on $\mathbb{R}^{d}$ by $h_{\omega^{(0)}}^{n}(y)=f\left(L_{n}\left(\omega^{(0)}\right), y\right)$. Then, noting that $f$ is bounded and Lipschitz continuous, Lemma 6.6 and (6.11) imply that

$$
\int h_{\omega^{(0)}}^{n}\left(L_{n}^{\prime}\left(\omega^{(1)}\right)\right) Q\left(\omega^{(0)}, \mathrm{d} \omega^{(1)}\right)-\int h_{\omega^{(0)}}^{n}\left(\left\|\phi_{2}\right\| \sqrt{\theta^{-1}} \widetilde{v}_{(\rho)}\left(\omega^{(0)}\right) y\right) \mathfrak{N}_{d}(\mathrm{~d} y) \rightarrow 0
$$

Since $f$ is bounded and $P\left(\Omega_{0}\right)=1$, this convergence and the bounded convergence theorem yield

$$
\Psi_{n}(1)-E\left[\int f\left(L_{n},\left\|\phi_{2}\right\| \sqrt{\theta^{-1}} \widetilde{v}_{(\rho)} y\right) \mathfrak{N}_{d}(\mathrm{~d} y) \mid \mathcal{G}_{\rho_{n}}\right] \rightarrow^{p} 0
$$

Next, since $f$ is Lipschitz continuous and $\widetilde{v}$ is cadlag and bounded, by (6.9) we obtain

$$
\begin{equation*}
\Psi_{n}(1)-E\left[\int f\left(L_{n},\left\|\phi_{2}\right\| \sqrt{\theta^{-1}} \widetilde{v}_{\rho_{n}} y\right) \mathfrak{N}_{d}(\mathrm{~d} y) \mid \mathcal{G}_{\rho_{n}}\right] \rightarrow^{p} 0 \tag{6.18}
\end{equation*}
$$

Now, noting that $W$ is a standard $d^{\prime}$-dimensional Brownian motion with respect to $\left(\mathcal{G}_{t}\right)$, by the strong Markov property of a Brownian motion $\left(W_{\rho_{n}+t}-W_{\rho_{n}}\right)_{t \geq 0}$ is independent of $\mathcal{G}_{\rho_{n}}$, hence
we have

$$
\begin{equation*}
E\left[\int f\left(L_{n},\left\|\phi_{2}\right\| \sqrt{\theta^{-1}} \widetilde{v}_{\rho_{n}} y\right) \mathfrak{N}_{d}(\mathrm{~d} y) \mid \mathcal{G}_{\rho_{n}}\right]=\int f\left(x,\left\|\phi_{2}\right\| \sqrt{\theta^{-1}} \widetilde{v}_{\rho_{n}} y\right) \mathbb{P}^{n}(\mathrm{~d} x) \mathfrak{N}_{d}(\mathrm{~d} y) \tag{6.19}
\end{equation*}
$$

where $\mathbb{P}^{n}$ is the law of $L_{n}$ under $P^{(0)}$. Then, again using the Lipschitz continuity of $f$ and the cadlag property of $\widetilde{v}$ as well as (6.9), we obtain

$$
\int f\left(x,\left\|\phi_{2}\right\| \sqrt{\theta^{-1}} \widetilde{v}_{\rho_{n}} y\right) \mathbb{P}^{n}(\mathrm{~d} x) \mathfrak{N}_{d}(\mathrm{~d} y)-\int f\left(x,\left\|\phi_{2}\right\| \sqrt{\theta^{-1}} \widetilde{v}_{(\rho)} y\right) \mathbb{P}^{n}(\mathrm{~d} x) \mathfrak{N}_{d}(\mathrm{~d} y) \rightarrow^{p} 0
$$

hence (6.12) yields

$$
\begin{align*}
& \int f\left(x,\left\|\phi_{2}\right\| \sqrt{\theta^{-1}} \widetilde{v}_{\rho_{n}} y\right) \mathbb{P}^{n}(\mathrm{~d} x) \mathfrak{N}_{d}(\mathrm{~d} y)  \tag{6.20}\\
& \quad \rightarrow^{p} \int f\left(\left\|\phi_{1}\right\| \sqrt{\theta G_{(\rho)}} x,\left\|\phi_{2}\right\| \sqrt{\theta^{-1}} \widetilde{v}_{(\rho)} y\right) \mathfrak{N}_{d^{\prime}}(\mathrm{d} x) \mathfrak{N}_{d}(\mathrm{~d} y)
\end{align*}
$$

Equations (6.18)-(6.20) imply that (6.17) holds true, and thus we complete the proof.

### 6.2.3. A joint convergence result and the proof of Proposition 6.2

In this subsection, we prove a joint convergence result for the pair $\left(n^{1 / 4} \mathbf{L}[M]^{n},\left(z_{r-}^{n}, z_{r-}^{\prime n}, z_{r+}^{n}\right.\right.$, $\left.z_{r+}^{\prime n}\right)_{r \geq 1}$ ) and complete the proof of Proposition 6.2.

For the proof, we use some elementary results on the Skorokhod topology. For any $k \in \mathbb{N}$, denote by $\mathbb{D}^{k}$ (resp., $\mathbb{D}^{k \times k}$ ) the space of $\mathbb{R}^{k}$-valued (resp., $\mathbb{R}^{k} \otimes \mathbb{R}^{k}$-valued) cadlag functions on $\mathbb{R}_{+}$equipped with the Skorokhod topology. For any $x \in \mathbb{D}^{k}$ and any $t \in \mathbb{R}_{+}$, we define the function $x^{t}$ by $x^{t}(s)=x(s \wedge t)$ for $s \in \mathbb{R}_{+}$. We evidently have $x^{t} \in \mathbb{D}^{k}$. On the other hand, for any $S \geq 0$ we define the function $\pi_{S}$ from $\mathbb{D}^{d \times d}$ into itself by $\pi_{S}(x)(t)=x(t-S) 1_{\{t \geq S\}}$.

The following two lemmas can be shown using basic properties of the Skorokhod topology, so we omit the proofs.

Lemma 6.8. The map $\mathbb{R}_{+} \times \mathbb{D}^{k} \ni(t, x) \mapsto x^{t} \in \mathbb{D}^{k}$ is continuous at every point $(t, x) \in \mathbb{R}_{+} \times \mathbb{D}^{k}$ such that $x$ is continuous at $t$, where the space $\mathbb{R}_{+} \times \mathbb{D}^{k}$ is equipped with the product topology.

Lemma 6.9. $\pi_{S}$ is a continuous function of $\mathbb{D}^{d \times d}$ into itself.
Now we are ready to prove the following joint convergence result:
Proposition 6.7. Suppose that [SA1]-[SA3] and [A4] are satisfied. Then

$$
\left(n^{1 / 4} \mathbf{L}[M]^{n},\left(z_{r-}^{n}, z_{r-}^{\prime n}, z_{r+}^{n}, z_{r+}^{\prime n}\right)_{r \geq 1}\right) \rightarrow^{d_{s}}\left(\mathcal{W},\left(z_{r-}, z_{r-}^{\prime}, z_{r+}, z_{r+}^{\prime}\right)_{r \geq 1}\right)
$$

as $n \rightarrow \infty$ for the product topology on the space $\mathbb{D}^{d \times d} \times\left(\mathbb{R}^{2\left(d+d^{\prime}\right)}\right)^{\mathbb{N}}$, where

$$
\begin{cases}z_{r-}=\psi_{2}^{-1} \sqrt{\Phi_{22} \theta G_{S_{r}-}} \Psi_{r-}, & z_{r+}=\psi_{2}^{-1} \sqrt{\Phi_{22} \theta G_{S_{r}}} \Psi_{r+} \\ z_{r-}^{\prime}=\psi_{2}^{-1} \sqrt{\Phi_{12} \theta^{-1}} \tilde{v}_{S_{r}-} \Psi_{r-}^{\prime}, & z_{r+}^{\prime}=\psi_{2}^{-1} \sqrt{\Phi_{12} \theta^{-1}} \widetilde{v}_{S_{r}} \Psi_{r+}^{\prime}\end{cases}
$$

Proof. Step 1. It suffices to prove

$$
\begin{equation*}
\left(n^{1 / 4} \mathbf{L}[M]^{n},\left(z_{r-}^{n}, z_{r-}^{\prime n}, z_{r+}^{n}, z_{r+}^{\prime n}\right)_{r \in \mathcal{R}}\right) \rightarrow^{d_{s}}\left(\mathcal{W},\left(z_{r-}, z_{r-}^{\prime}, z_{r+}, z_{r+}^{\prime}\right)_{r \in \mathcal{R}}\right) \tag{6.21}
\end{equation*}
$$

in $\mathbb{D}^{d \times d} \times \mathbb{R}^{2\left(d^{\prime}+d\right) \# \mathcal{R}}$ for any finite subset $\mathcal{R}$ of $\mathbb{N}$, and we prove this by induction on the number $\# \mathcal{R}$ of the elements in the set $\mathcal{R}$. First, (6.21) holds true when $\# \mathcal{R}=0$ due to Proposition 6.5. Next, let $J \in \mathbb{N}$ and assume that (6.21) holds true when $\# \mathcal{R}=J-1$. Then, we need to prove (6.21) for the case that $\# \mathcal{R}=J$. We write $\mathcal{R}=\left\{r_{1}, \ldots, r_{J}\right\}$ with $S_{r_{1}}<\cdots<S_{r_{J}}$.

Step 2. Before stating the detailed proof, we briefly explain the intuition behind the proof. The basic idea is the same as in the proof of Theorem 4.3 .1 from [31]. Namely, for each $\beta>0$ we set $S^{\beta-}=\left(S_{r_{J}}-\beta\right)_{+}$and $S^{\beta+}=S^{\beta+}$, and divide $n^{1 / 4} \mathbf{L}[M]^{n}$ into the summands containing the data observed in the interval $\left[S^{\beta-}, S^{\beta+}\right]$ and the remaining ones. Then we prove the negligibility of the former part (as $\beta \rightarrow 0$ ) and the joint limit theorem of the latter part and $\left(z_{r-}^{n}, z_{r-}^{\prime n}, z_{r+}^{n}, z_{r+}^{\prime n}\right)_{r \in \mathcal{R}}$. More formally, we set $\widehat{\mathbf{L}}(\beta)^{n}=n^{1 / 4}\left(\mathbf{L}[M]^{n}\right)^{S^{\beta+}}-n^{1 / 4}\left(\mathbf{L}[M]^{n}\right)^{S^{\beta-}}$ and $\widehat{\mathcal{W}}(\beta)=\mathcal{W}^{S^{\beta+}}-\mathcal{W}^{S^{\beta-}}$, and show that

$$
\begin{array}{r}
\limsup \limsup _{\beta \rightarrow 0} P\left(\sup _{n \rightarrow \infty}\left\|\widehat{\mathbf{L}}(\beta)_{t}^{n}\right\|>\eta\right)=0, \\
\limsup _{\beta \rightarrow 0} P\left(\sup _{0 \leq t \leq T}\left\|\widehat{\mathcal{W}}(\beta)_{t}\right\|>\eta\right)=0 \tag{6.22}
\end{array}
$$

for any $T, \eta>0$ and that

$$
\begin{gather*}
\left(n^{1 / 4} \mathbf{L}[M]^{n}-\widehat{\mathbf{L}}(\beta)^{n},\left(z_{r-}^{n}, z_{r-}^{\prime n}, z_{r+}^{n}, z_{r+}^{\prime n}\right)_{r \in \mathcal{R}}\right) \\
\quad \rightarrow^{d_{s}}\left(\mathcal{W}-\widehat{\mathcal{W}}(\beta),\left(z_{r-}, z_{r-}^{\prime}, z_{r+}, z_{r+}^{\prime}\right)_{r \in \mathcal{R}}\right) \tag{6.23}
\end{gather*}
$$

in $\mathbb{D}^{d \times d} \times \mathbb{R}^{2\left(d^{\prime}+d\right) J}$ as $n \rightarrow \infty$ for any fixed $\beta>0$. Then, Proposition 2.2 .4 of [31] yields (6.21).
However, to prove (6.23) we need a different approach from the one of [31] because we cannot argue conditionally on the increments of $W$ consisting of the observations in [ $S^{\beta-}, S^{\beta+}$ ] as [31] do, which is due to the time endogeneity. For this reason we further decompose $n^{1 / 4} \mathbf{L}[M]^{n}-$ $\widehat{\mathbf{L}}(\beta)^{n}$ as $n^{1 / 4} \mathbf{L}[M]^{n}-\widehat{\mathbf{L}}(\beta)^{n}=\check{\mathbf{L}}^{n}+\widetilde{\mathbf{L}}(\beta)^{n}$, where $\check{\mathbf{L}}^{n}=\left(n^{1 / 4} \mathbf{L}[M]^{n}\right)^{S^{\beta-}}$. Roughly speaking, $\check{\mathbf{L}}^{n}$ consists of the data observed before $S^{\beta-}$, while $\widetilde{\mathbf{L}}(\beta)^{n}$ consists of those observed after $S^{\beta+}$. By Proposition VI-1.23 of [32] and the continuous mapping theorem (6.23) follows once we show that

$$
\begin{equation*}
\left(\check{\mathbf{L}}^{n}, \tilde{\mathbf{L}}(\beta)^{n},\left(z_{r-}^{n}, z_{r-}^{\prime n}, z_{r+}^{n}, z_{r+}^{\prime n}\right)_{r \in \mathcal{R}}\right) \rightarrow^{d_{s}}\left(\mathcal{W}^{S_{r_{J}}}, \tilde{\mathcal{W}}(\beta),\left(z_{r-}, z_{r-}^{\prime}, z_{r+}, z_{r+}^{\prime}\right)_{r \in \mathcal{R}}\right) \tag{6.24}
\end{equation*}
$$

in $\mathbb{D}^{d \times d} \times \mathbb{D}^{d \times d} \times \mathbb{R}^{2\left(d^{\prime}+d\right) J}$ as $n \rightarrow \infty$. The strategy of the proof of (6.24) is, roughly speaking, as follows. We first prove a stable limit theorem for $\widetilde{\mathbf{L}}(\beta)^{n}$ conditionally on $\mathcal{F}_{S^{\beta+}}$ (this will be done in step 5; the assumption [A4](i) is necessary for this part). Then we obtain a joint stable limit theorem for $\widetilde{\mathbf{L}}(\beta)^{n}$ and $\left(z_{r_{J}-}^{n}, z_{r_{J}-}^{\prime n}, z_{r_{J}+}^{n}, z_{r_{J}+}^{\prime \prime}\right)$ conditionally on $\mathcal{G}_{S_{r_{J}}}$ by virtue of Lemma 6.7 (step 7). Finally, from the assumption of the induction we will obtain the desired result (step 8).

Step 3 . We begin with proving (6.22). The second equation immediately follows from the continuity of the process $\mathcal{W}$. On the other hand, for any $\beta>0$ we have $\left(S^{\beta-}, S^{\beta+}, n^{1 / 4} \mathbf{L}[M]^{n}\right) \rightarrow{ }^{d_{s}}$ $\left(S^{\beta-}, S^{\beta+}, \mathcal{W}\right)$ as $n \rightarrow \infty$ in $\mathbb{R}_{+} \times \mathbb{R}_{+} \times \mathbb{D}^{d \times d}$ by Proposition 6.5 , hence Lemma 6.8 and the continuous mapping theorem imply that $\left(n^{1 / 4}\left(\mathbf{L}[M]^{n}\right)^{S^{\beta+}}, n^{1 / 4}\left(\mathbf{L}[M]^{n}\right)^{S^{\beta-}}\right) \rightarrow^{d_{s}}$ $\left(\mathcal{W}^{S^{\beta+}}, \mathcal{W}^{S^{\beta-}}\right)$ as $n \rightarrow \infty$ in $\mathbb{D}^{d \times d} \times \mathbb{D}^{d \times d}$. Therefore, Propositions VI-1.23 and VI-2.4 of [32] as well as the continuous mapping theorem yield $\sup _{0 \leq t \leq T}\left\|\widehat{\mathbf{L}}(\beta)_{t}^{n}\right\| \rightarrow^{d_{s}} \sup _{0 \leq t \leq T}\left\|\mathcal{W}(\beta)_{t}\right\|$ as $n \rightarrow \infty$. In particular, we have

$$
\limsup _{n \rightarrow \infty} P\left(\sup _{0 \leq t \leq T}\left\|\widehat{\mathbf{L}}(\beta)_{t}^{n}\right\|>\eta\right) \leq P\left(\sup _{0 \leq t \leq T}\left\|\mathcal{W}(\beta)_{t}\right\| \geq \eta\right)
$$

hence we also obtain the first equation (6.22).
Step 4 . Now we start the proof of (6.24). First, due to the property of the product topology it suffices to prove the following convergence:

$$
\begin{equation*}
E\left[\zeta f_{1}\left(\check{\mathbf{L}}^{n}\right) f_{2}\left(\widetilde{\mathbf{L}}(\beta)^{n}\right) \prod_{j=1}^{J} Y_{j-}^{n} Y_{j+}^{n}\right] \rightarrow E\left[\zeta f_{1}\left(\mathcal{W}^{S_{r_{J}}}\right) f_{2}(\tilde{\mathcal{W}}(\beta)) \prod_{j=1}^{J} Y_{j-} Y_{j+}\right] \tag{6.25}
\end{equation*}
$$

as $n \rightarrow \infty$, where $\zeta$ is any bounded $\mathcal{F}^{(0)}$-measurable variable, $f_{1}$ and $f_{2}$ are bounded Lipschitz functions on $\mathbb{D}^{d \times d}$, and $Y_{j \pm}^{n}=F_{j \pm}\left(z_{r_{j} \pm}^{n}, z_{r_{j} \pm}^{\prime n}\right)$ and $Y_{j \pm}=F_{j \pm}\left(z_{r_{j} \pm}, z_{r_{j} \pm}^{\prime}\right)$ with $F_{j-}$ and $F_{j+}$ being bounded Lipschitz functions on $\mathbb{R}^{d^{\prime}+d}$ for every $j=1, \ldots, J$.

Step 5 . We begin with proving

$$
\begin{equation*}
E\left[\zeta f_{1}\left(\check{\mathbf{L}}^{n}\right) f_{2}\left(\widetilde{\mathbf{L}}(\beta)^{n}\right) \prod_{j=1}^{J} Y_{j-}^{n} Y_{j+}^{n}\right]-E\left[\zeta f_{1}\left(\check{\mathbf{L}}^{n}\right) f_{2}(\tilde{\mathcal{W}}(\beta)) \prod_{j=1}^{J} Y_{j-}^{n} Y_{j+}^{n}\right] \rightarrow 0 \tag{6.26}
\end{equation*}
$$

as $n \rightarrow \infty$. First, we introduce some notation. For any $d$-dimensional processes $U, V$, any $u, v \in$ $\left\{g, g^{\prime}\right\}$ and any $k, l=1, \ldots, d$, we define the process $\widetilde{\mathbb{L}}_{u, v}^{(k, l)}(U, V)^{n}$ in the same way as that of $\mathbb{L}_{u, v}^{(k, l)}(U, V)^{n}$ with replacing $\left(T_{p}\right)_{p \geq 0}$ by $\left(\widetilde{T}_{p}\right)_{p \geq 0}:=\left(T_{i\left(S^{\beta+}\right)^{n}+1+p}-S^{\beta+}\right)_{p \geq 0}$. Also, for any process $V$ we define the process $\stackrel{\circ}{V}$ by $\stackrel{\circ}{V}_{t}=V_{S^{\beta+}+t}-V_{S^{\beta+}}$, and define the $\mathbb{R}^{d} \otimes \mathbb{R}^{d}$-valued process $\widetilde{\mathbf{L}}^{n}$ by

$$
\widetilde{\mathbf{L}}^{n, k l}=n^{1 / 4}\left\{\widetilde{\mathbb{L}}_{g, g}^{(k, l)}(\stackrel{\circ}{M}, \stackrel{\circ}{M})^{n}+\widetilde{\mathbb{L}}_{g, g^{\prime}}^{(k, l)}(\stackrel{\circ}{M}, \stackrel{\circ}{\mathfrak{E}})^{n}+\widetilde{\mathbb{L}}_{g, g^{\prime}}^{(l, k)}(\stackrel{\circ}{M}, \stackrel{\circ}{\mathfrak{E}})^{n}+\widetilde{\mathbb{L}}_{g^{\prime}, g^{\prime}}^{(k, l)}(\stackrel{\circ}{\mathfrak{E}}, \stackrel{\circ}{\mathfrak{E}})^{n}\right\} .
$$

Then it can easily be seen that $\tilde{\mathbf{L}}(\beta)^{n}-\pi_{S^{\beta+}}\left(\widetilde{\mathbf{L}}^{n}\right) \xrightarrow{\text { ucp }} 0$ as $n \rightarrow \infty$. Therefore, by the Lipschitz continuity of $f_{2}$ as well as the boundedness of $\zeta, f_{1}$ and $Y_{j \pm}^{n}$ we have

$$
\begin{equation*}
E\left[\zeta f_{1}\left(\check{\mathbf{L}}^{n}\right) f_{2}\left(\widetilde{\mathbf{L}}(\beta)^{n}\right) \prod_{j=1}^{J} Y_{j-}^{n} Y_{j+}^{n}\right]-E\left[\zeta f_{1}\left(\check{\mathbf{L}}^{n}\right) f_{2}\left(\pi_{S^{\beta+}}\left(\widetilde{\mathbf{L}}^{n}\right)\right) \prod_{j=1}^{J} Y_{j-}^{n} Y_{j+}^{n}\right] \rightarrow 0 \tag{6.27}
\end{equation*}
$$

as $n \rightarrow \infty$.

Now we consider a regular conditional probability $p^{(0)}\left(\omega^{(0)}, \cdot\right)$ of $P^{(0)}$ given $\mathcal{F}_{S^{\beta+}}^{(0)}$. Such one exists because of the assumption [A4](i). We also consider a filtration $\left(\mathcal{F}_{t}^{(0)}\right)_{t \geq 0}$ of $\mathcal{F}^{(0)}$ defined by $\stackrel{\circ}{\mathcal{F}}_{t}^{(0)}=\mathcal{F}_{S^{\beta+}+t}^{(0)}$, and for each $\omega_{0} \in \Omega^{(0)}$ we introduce a stochastic basis $\mathcal{B}_{\omega_{0}}^{(0)}:=$ $\left(\Omega^{(0)}, \mathcal{F}^{(0)},\left(\stackrel{\circ}{\mathcal{F}}_{t}^{(0)}\right), p^{(0)}\left(\omega_{0}, \cdot\right)\right)$. For each $t \in \mathbb{R}_{+}$we also introduce a transition probability $\stackrel{\circ}{Q}_{t}\left(\omega^{(0)}, \mathrm{d} u\right)$ from $\left(\Omega^{(0)}, \circ_{f}^{(0)}\right)$ into $\mathbb{R}^{d}$ by setting $\stackrel{\circ}{Q}_{t}\left(\omega^{(0)}, A\right)=Q_{S^{\beta+}\left(\omega^{(0)}\right)+t}\left(\omega^{(0)}, A\right)$ for each Borel set $A$ of $\mathbb{R}^{d}$. Note that the process $\left(\mathscr{\circ}_{t}(\cdot, A)\right)_{t \geq 0}$ is $\left(\stackrel{\circ}{\mathcal{F}}_{t}^{(0)}\right)$-progressively measurable because of [A4](ii) and Theorem IV-57 of [41]. Now, by replacing $\mathcal{B}^{(0)}, Q_{t}\left(\omega^{(0)}, \mathrm{d} u\right)$ and $\left(\mathcal{T}_{i}^{n}\right)$ with $\mathcal{B}_{\omega_{0}}^{(0)}, \stackrel{\circ}{Q}_{t}\left(\omega^{(0)}, \mathrm{d} u\right)$ and the increasing reordering of $\left\{\tilde{\tau}_{p}^{k}:=\tau_{i\left(S^{\beta+}\right)^{n}+1+p}^{k}-S^{\beta+}: k=1, \ldots, d\right.$ and $p \geq 0\}$, respectively, we introduce the new stochastic basis $\mathcal{B}_{\omega_{0}}=\left(\Omega, \mathcal{F},\left(\stackrel{\circ}{\mathcal{F}}_{t}\right), P_{\omega_{0}}\right)$ instead of $\mathcal{B}$.

By the strong Markov property of a Brownian motion ${ }_{W}$ is a $d^{\prime}$-dimensional standard Brownian motion on $\mathcal{B}_{\omega_{0}}^{(0)}$. Moreover, defining the random measure $\stackrel{\circ}{\mu}$ by $\stackrel{\circ}{\mu}((0, t] \times A)=\mu\left(\left(S^{\beta+}, S^{\beta+}+\right.\right.$ $t] \times A$ ), the strong Markov property of a Poisson random measure implies that $\dot{\mu}$ is a Poisson random measure on $\mathcal{B}_{\omega_{0}}^{(0)}$ with compensator $v$. We also have

$$
\begin{aligned}
\stackrel{\circ}{M}_{t} & =\int_{0}^{t} \sigma_{S^{\beta+}+s} \mathrm{~d} \stackrel{\circ}{W}_{s}, \\
\sigma_{S^{\beta+}+t} & =\sigma_{S^{\beta+}}+\int_{0}^{t} \widetilde{b}_{S^{\beta+}+s} \mathrm{~d} s+\int_{0}^{t} \widetilde{\sigma}_{S^{\beta++s}} \mathrm{~d} \stackrel{\circ}{W}_{s}+\left(\stackrel{\circ}{\delta} 1_{\{|\check{\delta}| \leq 1\}}\right) \star(\stackrel{\circ}{\mu}-v)_{t}+\left(\stackrel{\tilde{\delta}}{( } 1_{\{|\tilde{\delta}|>1\}}\right) \star \stackrel{\circ}{\mu}_{t}, \\
G_{S^{\beta++t}} & =G_{S^{\beta+}}+\int_{0}^{t} \widehat{b}_{S^{\beta++s}} \mathrm{~d} s+\int_{0}^{t} \widehat{\sigma}_{S^{\beta+}+s} \mathrm{~d} \stackrel{\circ}{W}_{s}+\left(\stackrel{\circ}{\delta} 1_{\{|\dot{\delta}| \leq 1\}}\right) \star(\stackrel{\circ}{\mu}-v)_{t}+\left(\stackrel{\circ}{\delta} 1_{\{|\dot{\delta}|>1\}}\right) \star \stackrel{\circ}{\mu}_{t},
\end{aligned}
$$

where for a function $\eta$ on $\Omega^{(0)} \times \mathbb{R}_{+} \times E$ the function $\dot{\eta}$ on $\Omega^{(0)} \times \mathbb{R}_{+} \times E$ is defined by $\grave{\eta}\left(\omega^{(0)}, t, z\right)=\eta\left(\omega^{(0)}, S^{\beta+}\left(\omega^{(0)}\right)+t, z\right)$. Therefore, noting that for any $\mathcal{F}$-measurable variable $x$ and any sub $\sigma$-filed $\mathcal{H}$ of $\mathcal{F}$ we have $E_{P_{\omega_{0}}}[x \mid \mathcal{H}]=E[x \mid \mathcal{H}]$ as long as $\mathcal{F}_{S^{\beta+}}^{(0)} \subset \mathcal{H}$, it can easily been shown that the conditions [SA1]-[SA3] are satisfied with replacing $\mathcal{B}, X,\left(T_{p}\right),\left(\tau_{p}^{k}\right), G_{t}$ and $\chi_{t}$ by $\mathcal{B}_{\omega_{0}}, \stackrel{\circ}{M},\left(\widetilde{T}_{p}\right),\left(\widetilde{\tau}_{p}^{k}\right), G_{S^{\beta++t}}$ and $\chi_{S^{\beta++t}}$, respectively.

Consequently, Proposition 6.5 and Lemma 6.9 as well as the continuous mapping theorem yield

$$
E_{P_{\omega_{0}}}\left[\zeta f_{2}\left(\pi_{S_{r_{J}}\left(\omega_{0}\right)+\beta}\left(\widetilde{\mathbf{L}}^{n}\right)\right)\right] \rightarrow E_{p^{(0)}\left(\omega_{0}, \cdot\right)}\left[\zeta f_{2}\left(\pi_{S_{r_{J}}\left(\omega_{0}\right)+\beta}(\mathcal{W})\right)\right] .
$$

Therefore, noting that $E_{P_{\omega_{0}}}[\cdot]=E\left[\cdot \mid \mathcal{F}_{S^{\beta+}}\right]\left(\omega_{0}\right)$ and $E_{p^{(0)}\left(\omega_{0},\right)}\left[\zeta f_{2}\left(\pi_{S_{r_{J}}\left(\omega_{0}\right)+\beta}(\mathcal{W})\right)\right]=$ $E\left[\zeta f_{2}\left(\pi_{S_{J}\left(\omega_{0}\right)+\beta}(\mathcal{W})\right) \mid \mathcal{F}_{S^{\beta+}}\right]\left(\omega_{0}\right)$ for almost all $\omega_{0}$ (with respect to $P^{(0)}$ ) and that $f_{1}\left(\check{\mathbf{L}}^{n}\right) \prod_{j=1}^{J} Y_{j-}^{n} Y_{j+}^{n}$ is bounded and $\mathcal{F}_{S^{\beta+}}$-measurable (for sufficiently large $n$; note that $S^{\beta-}$ is an $\left(\mathcal{F}_{S^{\beta+}+t}\right)$-stopping time $)$, the bounded convergence theorem implies that

$$
E\left[\zeta f_{1}\left(\check{\mathbf{L}}^{n}\right) f_{2}\left(\pi_{S^{\beta+}}\left(\widetilde{\mathbf{L}}^{n}\right)\right) \prod_{j=1}^{J} Y_{j-}^{n} Y_{j+}^{n}\right]-E\left[\zeta f_{1}\left(\check{\mathbf{L}}^{n}\right) f_{2}\left(\pi_{S^{\beta+}}(\mathcal{W})\right) \prod_{j=1}^{J} Y_{j-}^{n} Y_{j+}^{n}\right] \rightarrow 0 .
$$

Since $\pi_{S^{\beta+}}(\mathcal{W})=\tilde{\mathcal{W}}(\beta)$, by (6.27) and the above convergence we obtain (6.26).
Step 6. In this step, we prove

$$
\begin{equation*}
E\left[\zeta^{\prime} Y_{J-}^{n} Y_{J+}^{n} \mid \mathcal{G}_{S_{r_{J}}^{\dagger}}\right] \rightarrow^{p} E\left[\zeta^{\prime} Y_{J-} Y_{J+} \mid \mathcal{G}_{S_{r_{J}}}^{(0)}\right] \tag{6.28}
\end{equation*}
$$

where $\zeta^{\prime}=E\left[\zeta f_{2}(\tilde{\mathcal{W}}(\beta)) \mid \mathcal{F}^{(0)}\right]$. Fix a constant $\alpha^{\prime}$ such that $1-\xi<\alpha^{\prime}<\left(\xi-\kappa-\frac{1}{2}\right) \wedge\left(\xi-\frac{3}{4}\right) \wedge$ $\varpi$. Then, set $\underline{k}_{n}=k_{n}-\left\lfloor n^{1 / 2-\alpha^{\prime}}\right\rfloor$ and define the $d^{\prime}$-dimensional variable $\underline{z}_{r_{J}-}^{n}=\left(\underline{z}_{r_{J}-}^{n, j}\right)_{1 \leq j \leq d^{\prime}}$ and the $d$-dimensional variable $\underline{z}_{r_{J}-}^{\prime n}=\left(\underline{z}_{J_{J}-}^{\prime n, k}\right)_{1 \leq k \leq d}$ by

$$
\underline{z}_{r^{\prime}-}^{n, j}=n^{1 / 4} \sum_{w=1}^{k_{n}}\left(\widetilde{\phi_{g, g}}\right)_{w}^{n} W^{j}\left(I_{i\left(S_{r}^{\dagger}\right)^{n}+w}\right), \quad \underline{z}_{r_{J}-}^{\prime n, k}=-\frac{n^{1 / 4}}{k_{n}} \sum_{w=1}^{\underline{k}_{n}}\left(\widetilde{\phi_{g^{\prime}, g}}\right)_{w}^{n} \epsilon_{\left.\tau_{i\left(S S_{r}^{\dagger}\right.}^{k}\right)^{n}+w}^{k}
$$

and put $\underline{Y}_{J-}^{n}=F_{J-}\left(\underline{z}_{r_{J}-}^{n}, \underline{z}_{r_{J}-}^{\prime n}\right)$. Since $E\left[\left\|z_{r_{J}-}^{n}-\underline{z}_{r_{J}-}^{n}\right\|^{2}\right] \lesssim n^{1-\alpha^{\prime}-\xi}$ and $E\left[\left\|z_{r_{J}-}^{\prime n}-\underline{z}_{r_{J}-}^{\prime n}\right\|^{2}\right] \lesssim$ $n^{-\alpha^{\prime}}$ by [SA1]-[SA3] and the optional sampling theorem, we have $Y_{J-}^{n}-\underline{Y}_{J-}^{n} \rightarrow^{p} 0$ due to the Lipschitz continuity of $F_{J-}$. Therefore, by virtue of the boundedness of $\zeta^{\prime}$ and $F_{J+}$, for the proof of (6.28) it is enough to prove

$$
\begin{equation*}
E\left[\zeta^{\prime} \underline{Y}_{J-}^{n} Y_{J+}^{n} \mid \mathcal{G}_{S_{r_{J}^{\prime}}^{\dagger}}\right] \rightarrow^{p} E\left[\zeta^{\prime} Y_{J-} Y_{J+} \mid \mathcal{G}_{S_{r_{J}-}}^{(0)}\right] . \tag{6.29}
\end{equation*}
$$

Now, Lemma 6.7 yields $E\left[\zeta^{\prime} Y_{J+}^{n} \mid \mathcal{G}_{S_{r_{J}}}\right] \rightarrow^{p} E\left[\zeta^{\prime} Y_{J+} \mid \mathcal{G}_{S_{r_{J}}}^{(0)}\right]$. Moreover, setting $\Omega_{n}=$ $\left\{T_{i\left(S_{r_{J}}^{\dagger}\right)+\bar{k}_{n}}<S_{r_{J}}\right\}$, we have

$$
E\left[\zeta^{\prime} \underline{Y}_{J-}^{n} Y_{J+}^{n} 1_{\Omega_{n}} \mid \mathcal{G}_{S_{r_{J}}^{\dagger}}\right]=E\left[\underline{Y}_{J-}^{n} E\left[\zeta^{\prime} Y_{J+}^{n} \mid \mathcal{G}_{S_{r_{J}}}\right] 1_{\Omega_{n}} \mid \mathcal{G}_{S_{r_{J}}^{\dagger}}\right]
$$

Since $\lim _{n} P\left(\Omega_{n}\right)=1$ by Lemma 6.4, the boundedness of $\zeta^{\prime}$ and $F_{J \pm}$ and the bounded convergence theorem imply that

$$
\begin{equation*}
E\left[\zeta^{\prime} \underline{Y}_{J-}^{n} Y_{J+}^{n} \mid \mathcal{G}_{S_{r_{J}^{\prime}}^{\dagger}}\right]-E\left[\zeta^{\prime \prime} \underline{Y}_{J-}^{n} \mid \mathcal{G}_{S_{r_{J}^{\prime}}^{\dagger}}\right] \rightarrow^{p} 0 \tag{6.30}
\end{equation*}
$$

where $\zeta^{\prime \prime}=E\left[\zeta^{\prime} Y_{r_{J}+} \mid \mathcal{G}_{S_{r_{J}}}^{(0)}\right]$. On the other hand, Lemma 6.7 again yields $E\left[\zeta^{\prime \prime} \underline{Y}_{J-}^{n} \mid \mathcal{G}_{S_{r_{J}}^{\dagger}}\right] \rightarrow p$ $E\left[\zeta^{\prime \prime} Y_{J-} \mid \mathcal{G}_{S_{r_{J}-}}^{(0)}\right.$. Since $Y_{J-}^{n}-\underline{Y}_{J-}^{n} \rightarrow^{p} 0$ and $\zeta^{\prime \prime}$ and $F_{J-}$ are bounded, the bounded convergence theorem again implies that

$$
\begin{equation*}
E\left[\zeta^{\prime \prime} \underline{Y}_{J-}^{n} \mid \mathcal{G}_{S_{r_{J}}^{\dagger}}\right] \rightarrow^{p} E\left[\zeta^{\prime \prime} Y_{J-} \mid \mathcal{G}_{S_{r_{J}}}^{(0)}\right] \tag{6.31}
\end{equation*}
$$

Equations (6.30) and (6.31) yield (6.29).
Step 7. Set $\Omega_{n}^{\prime}=\left\{S_{r_{J-1}}+k_{n} \bar{r}_{n}<S_{r_{J}}^{\dagger}\right\}$ if $J>1$ and $\Omega_{n}^{\prime}=\Omega$ otherwise. Then we have

$$
E\left[\zeta f_{1}\left(\check{\mathbf{L}}^{n}\right) f_{2}(\tilde{\mathcal{W}}(\beta)) \prod_{j=1}^{J} Y_{j-}^{n} Y_{j+}^{n} ; \Omega_{n}^{\prime}\right]=E\left[f_{1}\left(\check{\mathbf{L}}^{n}\right) \prod_{j=1}^{J-1} Y_{j-}^{n} Y_{j+}^{n} E\left[\zeta^{\prime} Y_{J-}^{n} Y_{J+}^{n} \mid \mathcal{G}_{S_{J}^{\dagger}}\right] ; \Omega_{n}^{\prime}\right]
$$

Therefore, by (6.28) and the boundedness of $\zeta^{\prime}, f_{1}$ and $F_{j \pm}$ we obtain

$$
E\left[\zeta f_{1}\left(\check{\mathbf{L}}^{n}\right) f_{2}(\widetilde{\mathcal{W}}(\beta)) \prod_{j=1}^{J} Y_{j-}^{n} Y_{j+}^{n} ; \Omega_{n}^{\prime}\right]-E\left[\check{\zeta} f_{1}\left(\check{\mathbf{L}}^{n}\right) \prod_{j=1}^{J-1} Y_{j-}^{n} Y_{j+}^{n} ; \Omega_{n}^{\prime}\right] \rightarrow 0,
$$

where $\check{\zeta}=E\left[\zeta^{\prime} Y_{J-} Y_{J+} \mid \mathcal{G}_{S_{r_{J}-}}^{(0)}\right.$. Since $\lim _{n} P\left(\Omega_{n}^{\prime}\right)=1$, we conclude that

$$
\begin{equation*}
E\left[\zeta f_{1}\left(\check{\mathbf{L}}^{n}\right) f_{2}(\tilde{\mathcal{W}}(\beta)) \prod_{j=1}^{J} Y_{j-}^{n} Y_{j+}^{n}\right]-E\left[\check{\zeta} f_{1}\left(\check{\mathbf{L}}^{n}\right) \prod_{j=1}^{J-1} Y_{j-}^{n} Y_{j+}^{n}\right] \rightarrow 0 \tag{6.32}
\end{equation*}
$$

Step 8. Now we are ready to prove (6.25). From (6.26) and (6.32) it remains to prove

$$
\begin{equation*}
E\left[\check{\zeta} f_{1}\left(\check{\mathbf{L}}^{n}\right) \prod_{j=1}^{J-1} Y_{j-}^{n} Y_{j+}^{n}\right] \rightarrow E\left[\zeta f_{1}\left(\mathcal{W}^{S_{r_{J}}}\right) f_{2}(\widetilde{\mathcal{W}}(\beta)) \prod_{j=1}^{J} Y_{j-} Y_{j+}\right] \tag{6.33}
\end{equation*}
$$

By the assumption of the induction, we have

$$
\left(n^{1 / 4} \mathbf{L}[M]^{n},\left(z_{r_{j}-}^{n}, z_{r_{j}-}^{\prime n}, z_{r_{j}+}^{n}, z_{r_{j}+}^{\prime n}\right)_{j=1}^{J-1}\right) \rightarrow^{d_{s}}\left(\mathcal{W},\left(z_{r_{j}-}, z_{r_{j}-}^{\prime}, z_{r_{j}+}, z_{r_{j}+}^{\prime}\right)_{j=1}^{J-1}\right)
$$

in $\mathbb{D}^{d \times d} \times \mathbb{R}^{2\left(d^{\prime}+d\right)(J-1)}$ as $n \rightarrow \infty$. Therefore, by Lemma 6.8 and the continuous mapping theorem we obtain

$$
\left(\check{\mathbf{L}}^{n},\left(z_{r_{j}-}^{n}, z_{r_{j}-}^{\prime n}, z_{r_{j}+}^{n}, z_{r_{j}+}^{\prime n}\right)_{j=1}^{J-1}\right) \rightarrow^{d_{s}}\left(\mathcal{W}^{S_{r_{J}}},\left(z_{r_{j}-}, z_{r_{j}-}^{\prime}, z_{r_{j}+}, z_{r_{j}+}^{\prime}\right)_{j=1}^{J-1}\right)
$$

in $\mathbb{D}^{d \times d} \times \mathbb{R}^{2\left(d^{\prime}+d\right)(J-1)}$. This implies that $E\left[\check{\zeta} f_{1}\left(\check{\mathbf{L}}^{n}\right) \prod_{j=1}^{J-1} Y_{j-}^{n} Y_{j+}^{n}\right] \rightarrow E\left[\check{\zeta} f_{1}\left(\mathcal{W}^{S_{r_{J}}}\right) \times\right.$ $\prod_{j=1}^{J-1} Y_{j-} Y_{j+}$ ]. Now, since $\zeta f_{2}(\tilde{\mathcal{W}}(\beta))$ is independent of $Y_{J \pm}$ by construction, we have $\check{\zeta}=E\left[\zeta f_{2}(\tilde{\mathcal{W}}(\beta)) Y_{J-} Y_{J+} \mid \mathcal{G}_{S_{r_{J}-}}^{(0)}\right]$. Moreover, since $\zeta f_{2}(\tilde{\mathcal{W}}(\beta)) Y_{J-} Y_{J+}$ is independent of $Y_{1 \pm}, \ldots, Y_{(J-1) \pm}$ by construction, we conclude that

$$
E\left[\check{\zeta} f_{1}\left(\mathcal{W}^{S_{r_{J}}}\right) \prod_{j=1}^{J-1} Y_{j-} Y_{j+}\right]=E\left[\zeta f_{1}\left(\mathcal{W}^{S_{r_{J}}}\right) f_{2}(\tilde{\mathcal{W}}(\beta)) \prod_{j=1}^{J} Y_{j-} Y_{j+}\right]
$$

hence we obtain (6.33).
Proof of Proposition 6.2. First, by Propositions 6.4-6.7 as well as properties of stable convergence, we have

$$
\begin{aligned}
& \left(n^{1 / 4}\left(\Xi[C(m)]^{n}-[M, M]-\frac{\psi_{1}}{\psi_{2} k_{n}^{2}}[Y, Y]^{n}\right),\left(\eta_{-}(n, r), \eta_{-}^{\prime}(n, r), \eta_{+}(n, r), \eta_{+}^{\prime}(n, r)\right)_{r \geq 1}\right) \\
& \quad \rightarrow^{d_{s}}\left(\mathcal{W},\left(\sigma_{S_{r}-}-z_{r-}^{n}, z_{r-}^{\prime n}, \sigma_{S_{r}} z_{r+}^{n}, z_{r+}^{\prime n}\right)_{r \geq 1}\right) .
\end{aligned}
$$

The second claim immediately follows from this convergence. On the other hand, since equations (6.5)-(6.6) holds on the set $\Omega_{n}(t, m)$ (recall that $\Omega_{n}(t, m)$ is defined at the beginning of Section 6.2.1), we obtain

$$
\begin{align*}
& n^{1 / 4}\left(\Xi[X(m)]_{t}^{n}-\left(\Xi_{g, g}^{(k, l)}(J(m), J(m))_{t}^{n}\right)_{1 \leq k, l, \leq d}-[M, M]_{t}-\frac{\psi_{1}}{\psi_{2} k_{n}^{2}}[Y, Y]_{t}^{n}\right)  \tag{6.34}\\
& \quad \rightarrow^{d_{s}} \mathcal{W}_{t}+\mathcal{Z}(m)_{t}
\end{align*}
$$

by the continuous mapping theorem and the fact that (3.2) yields $\lim _{n \rightarrow \infty} P\left(\Omega_{n}(t, m)\right)=1$.
Next, let $\Omega_{n}^{\prime}(m)$ be the set on which $\left|S_{r_{1}}-S_{r_{2}}\right|>k_{n} \bar{r}_{n}$ for any $r_{1}, r_{2} \in \mathcal{R}_{m}$ such that $r_{1} \neq r_{2}$ and $S_{r_{1}}, S_{r_{2}}<\infty$. Then we have

$$
\Xi_{g, g}^{(k, l)}(J(m), J(m))_{t}^{n}=\frac{1}{\psi_{2}} \sum_{r \in \mathcal{R}_{m}: S_{r} \leq t}\left(\frac{1}{k_{n}} \sum_{i=1}^{k_{n}-1} g_{i}^{n} g_{i}^{n}\right) \Delta X_{S_{r}}^{k} \Delta X_{S_{r}}^{l}
$$

on the set $\Omega_{n}^{\prime}(m) \cap \Omega_{n}(t, m)$. Since $\frac{1}{k_{n}} \sum_{i=1}^{k_{n}-1} g_{i}^{n} g_{i}^{n}=\psi_{2}+\mathrm{O}\left(k_{n}^{-1}\right)$ by the Lipschitz continuity of $g$ and $\lim _{n} P\left(\Omega_{n}^{\prime}(m) \cap \Omega_{n}(t, m)\right)=1$, we obtain

$$
\begin{equation*}
n^{1 / 4}\left\{\Xi_{g, g}^{(k, l)}(J(m), J(m))_{t}^{n}-\sum_{r \in \mathcal{R}_{m}: S_{r} \leq t} \Delta X_{S_{r}}^{k} \Delta X_{S_{r}}^{l}\right\} \rightarrow^{p} 0 \tag{6.35}
\end{equation*}
$$

Finally, since

$$
\left[X(m)^{k}, X(m)^{l}\right]_{t}=\left[M^{k}, M^{l}\right]_{t}+\sum_{r \in \mathcal{R}_{m}: S_{r} \leq t} \Delta X_{S_{r}}^{k} \Delta X_{S_{r}}^{l},
$$

(6.34) and (6.35) imply the first claim of the proposition.

### 6.3. Proof of Proposition 6.3

We decompose the target quantity as

$$
\begin{align*}
& \boldsymbol{\Xi}[X]_{t}^{n, k l}-\boldsymbol{\Xi}[X(m)]_{t}^{n, k l} \\
& \quad=\left\{\boldsymbol{\Xi}_{g, g}^{(k, l)}(X, X)^{n}-\Xi_{g, g}^{(k, l)}(X(m), X(m))^{n}\right\}+\boldsymbol{\Xi}_{g, g^{\prime}}^{(k, l)}(Z(m), \mathfrak{E})^{n}+\Xi_{g, g^{\prime}}^{(l, k)}(Z(m), \mathfrak{E})^{n} . \tag{6.36}
\end{align*}
$$

We start by proving the negligibility of the second and the third terms in the right-hand side of the above equation, which can be shown by an easy calculation.

Lemma 6.10. Under the assumptions of Proposition 6.3, it holds that

$$
\limsup _{m \rightarrow \infty} \limsup _{n \rightarrow \infty} P\left(n^{1 / 4}\left|\Xi_{g, g^{\prime}}^{(k, l)}(Z(m), \mathfrak{E})_{t}^{n}\right|>\eta\right)=0
$$

for any $t, \eta>0$.

Proof. First, since $E_{0}\left[\overline{\mathfrak{E}}\left(g^{\prime}\right)_{i}^{l} \overline{\mathfrak{E}}\left(g^{\prime}\right)_{j}^{l}\right]=0$ if $|i-j| \geq k_{n}$ and $\left|E_{0}\left[\overline{\mathfrak{E}}\left(g^{\prime}\right)_{i}^{l} \overline{\mathfrak{E}}\left(g^{\prime}\right)_{j}^{l}\right]\right| \lesssim k_{n}^{-1}$ by [SA3] and the definition of $\mathfrak{E}$, we have

$$
E\left[\left|n^{1 / 4} \Xi_{g, g^{\prime}}^{(k, l)}(Z(m), \mathfrak{E})_{t}^{n}\right|^{2}\right] \lesssim \frac{\sqrt{n}}{k_{n}^{2}} E\left[\sum_{i=1}^{N_{t}^{n}-k_{n}+1}\left|\overline{Z(m)}(g)_{i}^{k}\right|^{2}\right] .
$$

Next, the definition of $Z(m)$ and the optimal sampling theorem yield

$$
\begin{aligned}
E\left[\frac{\sqrt{n}}{k_{n}^{2}} \sum_{i=1}^{N_{t}^{n}-k_{n}+1}\left|\overline{Z(m)}(g)_{i}^{k}\right|^{2}\right] & \leq E\left[\frac{\sqrt{n}}{k_{n}^{2}} \sum_{i=1}^{\infty}\left|\sum_{p=0}^{k_{n}-1} g_{p}^{n} Z(m)^{k}\left(I_{i+p}(t)\right)\right|^{2}\right] \\
& \leq \frac{\sqrt{n}}{k_{n}^{2}}\|g\|_{\infty} \bar{\gamma}_{m} E\left[\sum_{i=1}^{\infty} \sum_{p=0}^{k_{n}-1}\left|I_{i+p}(t)\right|\right] \leq \frac{\sqrt{n}}{k_{n}} t\|g\|_{\infty} \bar{\gamma}_{m},
\end{aligned}
$$

where $\bar{\gamma}_{m}=\int_{A_{m}^{c}} \gamma(z)^{2} \lambda(\mathrm{~d} z)$. Since $\sqrt{n} / k_{n}=\mathrm{O}(1)$ as $n \rightarrow \infty$ and $\lim _{m} \bar{\gamma}_{m}=0$ by the dominated convergence theorem, we conclude that

$$
\limsup _{m \rightarrow \infty} \limsup _{n \rightarrow \infty} E\left[\frac{\sqrt{n}}{k_{n}^{2}} \sum_{i=1}^{N_{t}^{n}-k_{n}+1}\left|\overline{Z(m)}(g)_{i}^{k}\right|^{2}\right]=0 .
$$

Therefore, the Chebyshev inequality implies the desired result.
Next, we prove the negligibility of the term $\Xi_{g, g}^{(k, l)}(X, X)_{t}^{n}-\Xi_{g, g}^{(k, l)}(X(m), X(m))_{t}^{n}$. We further decompose it as

$$
\begin{align*}
& \Xi_{g, g}^{(k, l)}(X, X)_{t}^{n}-\Xi_{g, g}^{(k, l)}(X(m), X(m))_{t}^{n} \\
& \quad=\Xi_{g, g}^{(k, l)}(Z(m), X)_{t}^{n}+\Xi_{g, g}^{(k, l)}(X, Z(m))_{t}^{n}-\Xi_{g, g}^{(k, l)}(Z(m), Z(m))_{t}^{n} . \tag{6.37}
\end{align*}
$$

Therefore, using the decomposition

$$
\begin{equation*}
X_{t}=X_{0}+B_{t}^{\prime}+M_{t}+Z_{t}, \tag{6.38}
\end{equation*}
$$

where $B_{t}^{\prime}=\int_{0}^{t} b_{s}^{\prime} \mathrm{d} s, b_{s}^{\prime}=b_{s}+\int_{\{\|\delta(s, z)\|>1\}} \delta(s, z) \lambda(\mathrm{d} z)$ and $Z_{t}=\delta \star(\mu-v)_{t}$, it is enough to prove the negligibility of $\Xi_{g, g}^{(k, l)}(Z(m), V)_{t}^{n}$ for $V \in\left\{B^{\prime}, M, Z, Z(m)\right\}$. In the following we fix $V \in\left\{B^{\prime}, M, Z, Z(m)\right\}$.

Lemma 6.11. Assume [SA2]. Then,
(a) $\sup _{0 \leq h \leq h_{0}}\left\|V_{t}-V_{(t-h)_{+}}\right\|=\mathrm{O}_{p}\left(\sqrt{h_{0}}\right)$ as $h_{0} \downarrow 0$,
(b) $\sup _{0 \leq h \leq h_{0}}\left|\left[Z(m)^{k}, V^{l}\right]_{t}-\left[Z(m)^{k}, V^{l}\right]_{(t-h)_{+}}\right|=\mathrm{O}_{p}\left(h_{0}\right)$ as $h_{0} \downarrow 0$ for all $k, l$.

Proof. The claim is evident if $V=B^{\prime}$, so we assume that $V \neq B^{\prime}$. Then the Doob inequality and [SA2] yield $E\left[\sup _{0 \leq h \leq h_{0}}\left\|V_{t}-V_{(t-h)_{+}}\right\|^{2}\right] \lesssim h_{0}$, which implies (a). On the other hand, the Kunita-Watanabe and Schwarz inequalities as well as [SA2] yield $E\left[\sup _{0 \leq h \leq h_{0}} \mid\left[Z(m)^{k}, V^{l}\right]_{t}-\right.$ $\left.\left[Z(m)^{k}, V^{l}\right]_{(t-h)_{+}} \mid\right] \lesssim h_{0}$, which implies (b).

Lemma 6.12. $\sup _{1 \leq q \leq N_{t}^{n}+1}\left|C_{g, g}^{n}(V)_{q}^{k}\right|=\mathrm{O}_{p}(1)$ as $n \rightarrow \infty$ for any $t>0$ and $k=1, \ldots, d$.
Proof. This can be shown in the same manner as the proof of Lemma 6.8 from [36].
Lemma 6.13. Under the assumptions of Proposition 6.3, it holds that

$$
n^{1 / 4}\left\{\Xi_{g, g}^{(k, l)}(Z(m), V)_{t}^{n}-\mathbb{L}_{g, g}^{(k, l)}(Z(m), V)_{t}^{n}-\left[Z(m)^{k}, V^{l}\right]_{t}\right\} \rightarrow^{p} 0
$$

as $n \rightarrow \infty$ for any $t>0$.
Proof. Simple calculations and Lemma 6.11(a) yield

$$
\mathbf{B}:=\Xi_{g, g}^{(k, l)}(Z(m), V)_{t}^{n}-\mathbb{L}_{g, g}^{(k, l)}(Z(m), V)_{t}^{n}=\sum_{p=1}^{N_{t}^{n}+1} c_{g, g}^{n}(p, p) Z(m)^{k}\left(I_{p}\right)_{t} V^{l}\left(I_{p}\right)+\mathrm{o}_{p}\left(n^{-1 / 4}\right)
$$

Therefore, using Lemma 6.11(b), we can prove $\mathbf{B}=\left[Z(m)^{k}, V^{l}\right]_{t}+\mathrm{o}_{p}\left(n^{-1 / 4}\right)$ analogously to the proof of equation (6.24) from [36].

Lemma 6.14. Under the assumptions of Proposition 6.3, it holds that

$$
\begin{equation*}
\limsup _{m \rightarrow \infty} \limsup _{n \rightarrow \infty} P\left(n^{1 / 4}\left|\mathbb{L}_{g, g}^{(k, l)}(Z(m), V)_{t}^{n}\right|>\eta\right)=0 \tag{6.39}
\end{equation*}
$$

for any $t, \eta>0$.
Proof. First, if $V=B^{\prime}$, we can adopt an analogous argument to the proof of Lemma 6.10 from [36] and deduce $n^{1 / 4} \mathbb{L}_{g, g}^{(k, l)}(Z(m), V)_{t}^{n} \rightarrow^{p} 0$ as $n \rightarrow \infty$ for every $m$, so (6.39) holds true.

Next, we suppose that $V \neq B^{\prime}$. It is enough to prove

$$
\begin{align*}
& \limsup \limsup  \tag{6.40}\\
& m \rightarrow \infty  \tag{6.41}\\
& n \rightarrow \infty \\
& \underset{m \rightarrow \infty}{ } P\left(n^{1 / 4}\left|\mathbb{M}_{g, g}^{(k, l)}(Z(m), V)_{t}^{n}\right|>\eta^{\prime}\right)=0 \\
& \limsup \\
& \limsup \\
& n \rightarrow \infty \\
&
\end{align*}\left(n^{1 / 4}\left|\mathbb{M}_{g, g}^{(k, l)}(V, Z(m))_{t}^{n}\right|>\eta^{\prime}\right)=0,
$$

for any $\eta^{\prime}>0$. Since we can prove (6.41) in a similar manner to the proof of (6.40), we only prove (6.40).

Since $V$ is an $\left(\mathcal{F}_{t}^{(0)}\right)$-martingale for any $n$ due to [SA2], by the Lenglart inequality it suffices to show that

$$
\begin{equation*}
\limsup _{m \rightarrow \infty} \limsup _{n \rightarrow \infty} P\left(U(n, m)_{t}>\eta^{\prime}\right)=0 \tag{6.42}
\end{equation*}
$$

for any $\eta^{\prime}>0$, where $U(n, m)_{t}=\sqrt{n} \sum_{q=2}^{N_{t}^{n}+1} E\left[\left|C_{g, g}^{n}(Z(m))_{q}^{k} V^{l}\left(I_{q}\right)\right|^{2} \mid \mathcal{F}_{T_{q-1}}^{(0)}\right]$. To prove (6.42), for each $j \geq 1$ we set $\Lambda(j)_{q}^{n}=\left\{E\left[n \mid I_{q} \| \mathcal{F}_{T_{q-1}}^{(0)}\right] \leq j\right\}$ and decompose $U(n, m)$ as

$$
\begin{aligned}
U(n, m)_{t} & =\sqrt{n} \sum_{q=2}^{N_{t}^{n}+1} E\left[\left|C_{g, g}^{n}(Z(m))_{q}^{k} V^{l}\left(I_{q}\right)\right|^{2} \mid \mathcal{F}_{T_{q-1}}^{(0)}\right]\left(1_{\Lambda(j)_{q}^{n}}+1_{\left(\Lambda(j)_{q}^{n}\right) c}\right) \\
& =U(n, m, j)_{t}+U^{\prime}(n, m, j)_{t}
\end{aligned}
$$

First, we prove

$$
\begin{equation*}
\limsup _{m \rightarrow \infty} \limsup _{n \rightarrow \infty} P\left(U(n, m, j)_{t}>\eta^{\prime}\right)=0 \tag{6.43}
\end{equation*}
$$

for any fixed $j$ and any $\eta^{\prime}>0$. We have

$$
\begin{aligned}
E\left[U(n, m, j)_{t}\right] & =\sqrt{n} E\left[\sum_{q=2}^{N_{t}^{n}+1}\left|C_{g, g}^{n}(Z(m))_{q}^{k}\right|^{2} E\left[\left\langle V^{l}\right\rangle\left(I_{q}\right) \mid \mathcal{F}_{T_{p-1}}^{(0)}\right] 1_{\Lambda(j)_{q}^{n}}\right] \\
& \lesssim \frac{j}{\sqrt{n}} E\left[\sum_{q=2}^{\infty} \sum_{p=\left(q-k_{n}\right) \vee 1}^{q-1}\left|c_{u, v}^{n}(p, q)\right|^{2}\left\langle Z(m)^{k}\right\rangle\left(I_{p}(t)\right)\right] \\
& \lesssim \bar{\gamma}_{m} \frac{1}{\sqrt{n}} E\left[\sum_{q=2}^{\infty} \sum_{p=\left(q-k_{n}\right) \vee 1}^{q-1}\left|I_{p}(t)\right|\right] \leq \bar{\gamma}_{m} k_{n} n^{-1 / 2} t,
\end{aligned}
$$

hence it holds that $\limsup \operatorname{sim}_{m} \sup _{n} E\left[U(n, m, j)_{t}\right]=0$. Therefore, we obtain (6.43) by the Chebyshev inequality.

Next, we prove

$$
\begin{equation*}
\underset{j \rightarrow \infty}{\limsup } \limsup _{m \rightarrow \infty} \limsup _{n \rightarrow \infty} P\left(U^{\prime}(n, m, j)_{t}>\eta^{\prime}\right)=0 \tag{6.44}
\end{equation*}
$$

for any $\eta^{\prime}>0$. Since $E\left[n\left|I_{q}\right| \mid \mathcal{F}_{T_{q-1}}^{(0)}\right]=E\left[E\left[n\left|I_{q}\right| \mid \mathcal{G}_{T_{q-1}}^{(0)}\right] \mid \mathcal{F}_{T_{q-1}}^{(0)}\right]$, we have $\Lambda(j)_{q}^{n} \supset$ $\left\{E\left[n \mid I_{q} \| \mathcal{G}_{T_{q-1}}^{(0)}\right] \leq j\right\}$. Therefore,

$$
\begin{aligned}
U^{\prime}(n, m, j)_{t} & \leq \sqrt{n} \sum_{q=2}^{N_{t}^{n}+1} E\left[\left|C_{g, g}^{n}(Z(m))_{q}^{k} V^{l}\left(I_{q}\right)\right|^{2} \mid \mathcal{F}_{T_{q-1}}^{(0)}\right]\left(1_{\left\{G_{T_{q-1}}^{n}>j\right\}}+1_{\left\{q-1 \in \mathcal{N}^{n}\right\}}\right) \\
& =: U^{\prime}(n, m, j)_{t}^{(1)}+U^{\prime}(n, m, j)_{t}^{(2)}
\end{aligned}
$$

Since $\left\{U^{\prime}(n, m, j)^{(1)}>0\right\} \subset\left\{\sup _{0 \leq s \leq t} G_{s}^{n}>j\right\}$, we have $\lim \sup _{j} \lim \sup _{m} \lim \sup _{n} P\left(U^{\prime}(n, m\right.$, $\left.j)_{t}^{(1)}>0\right)=0$. On the other hand, [SA2] and (6.2) imply that

$$
U^{\prime}(n, m, j)_{t}^{(2)} \leq \sqrt{n} \bar{r}_{n} \sup _{1 \leq q \leq N_{t}^{n}+1}\left|C_{g, g}^{n}(Z(m))_{q}^{k}\right|^{2} \#\left(\mathcal{N}^{n} \cap\left\{q: T_{q} \leq t\right\}\right),
$$

hence Lemma 6.12, (6.1) and [A1](i) yield $\limsup _{n} P\left(U^{\prime}(n, m, j)_{t}^{(2)}>\eta^{\prime}\right)=0$. Consequently, we obtain (6.44).

From (6.43) we have $\limsup _{m} \limsup _{n} P\left(U(n, m)_{t}>\eta^{\prime}\right) \leq \limsup \operatorname{simsup}_{n} P\left(U^{\prime}(n\right.$, $m, j)_{t}>\eta^{\prime}$ ) for any $j \geq 1$ and any $\eta^{\prime}>0$. Hence, (6.44) yields (6.42), which completes the proof.

Proof of Proposition 6.3. (6.4) immediately follows from equations (6.38) and (6.36)-(6.37) as well as Lemmas 6.10 and 6.13-6.14. Equation (6.3) follows from the equation $\widetilde{E}\left[\mid \mathcal{Z}(m)_{t}^{k l}-\right.$ $\left.\left.\mathcal{Z}_{t}^{k l}\right|^{2} \mid \mathcal{F}^{(0)}\right]=\frac{1}{\psi_{2}^{2}} \sum_{r \notin \mathcal{R}_{m}: S_{r} \leq t}\left(\mathfrak{J}_{S_{r}}^{k k}+2 \mathfrak{J}_{S_{r}}^{k l}+\mathfrak{J}_{S_{r}}^{l l}\right)$ and the fact that $\sum_{r \notin \mathcal{R}_{m}: S_{r} \leq t}\left\|\Delta X_{S_{r}}\right\|^{2} \rightarrow^{p} 0$ as $m \rightarrow \infty$.

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[^0]:    *Note. We report the sample mean, standard deviation (SD) as well as the $95 \%$ and $99 \%$ coverages of the standardized statistics (5.2) included in the simulation study.

