# A TEST FOR THE RANK OF THE VOLATILITY PROCESS: THE RANDOM PERTURBATION APPROACH 

By Jean Jacod and Mark Podolskij<br>Institut de Mathématiques de Jussieu and Heidelberg University


#### Abstract

In this paper, we present a test for the maximal rank of the matrix-valued volatility process in the continuous Itô semimartingale framework. Our idea is based upon a random perturbation of the original high frequency observations of an Itô semimartingale, which opens the way for rank testing. We develop the complete limit theory for the test statistic and apply it to various null and alternative hypotheses. Finally, we demonstrate a homoscedasticity test for the rank process.


1. Introduction. In recent years asymptotic statistics for high frequency observations has received a lot of attention in the literature. This interest was mainly motivated by financial applications, where observations of stocks or currencies are available at very high frequencies. As under the no-arbitrage condition, prices processes must be semimartingales (see, e.g., [4]), a lot of research has been devoted to statistics of high frequency data of semimartingales. We refer to a recent book [10] for a comprehensive study of infill asymptotic for semimartingales.

This paper is devoted to testing for the maximal rank of the matrix-valued volatility process in the continuous Itô semimartingale framework, and more specifically for a $d$-dimensional continuous Itô semimartingale $X$ which is observed at equidistant times over a fixed time interval [ $0, T$ ]: we observe $\left(X_{i \Delta_{n}}\right)_{0 \leq i \leq\left[T / \Delta_{n}\right]}$, and the high-frequency approach consists in assuming $\Delta_{n} \rightarrow 0$.

A continuous Itô semimartingale can be written as

$$
\begin{equation*}
d X_{t}=b_{t} d t+\sigma_{t} d W_{t} \tag{1.1}
\end{equation*}
$$

where $W$ is a Brownian motion, and there are many representations of this form, with different Brownian motions $W$ and, accordingly, different volatility processes $\sigma$. What is "intrinsic" is the drift coefficient $b_{t}$ and the diffusion coefficient ("squared volatility") $c_{t}=\sigma_{t} \sigma_{t}^{*}$, in the sense that they are uniquely determined by $X$, up to a Lebesgue-null set of times (throughout the paper $\sigma_{t}^{*}$ denotes the transpose of the matrix $\sigma$ ).

[^0]For modeling purposes and economical interpretation, we would like to find, and often choose, the smallest possible dimension of the Brownian motion $W$ in the representation (1.1). Assuming further that $t \mapsto c_{t}$ is continuous, this smallest possible dimension is the supremum in time of the rank of the $\mathbb{R}^{d \times d}$-valued process $c$ over the time interval $[0, T]$. Although the estimation of the maximal rank of $c$, which is an integer, does not make much sense (see the discussion in Section 3.3), testing for the maximal rank is a well formulated statistical problem. Besides the obvious application of the rank test, we will link our testing procedure to other statistical problems, such as testing for market (in)completeness in financial mathematics, testing for integrated diffusions and testing the local volatility vs. stochastic volatility hypothesis (cf. Remark 3.9). We are further interested in homoscedasticity testing for the rank process.

A partial answer to this question was given in [9]. The authors of this paper studied the problem of testing the null hypothesis $\sup _{t \in[0, T)} \operatorname{rank}\left(c_{t}\right) \geq r_{0}$ against $\sup _{t \in[0, T)} \operatorname{rank}\left(c_{t}\right)<r_{0}$ for a given number $r_{0}$. However, their method does not extend to testing null hypotheses of other types, for example, $\sup _{t \in[0, T)} \operatorname{rank}\left(c_{t}\right)=r_{0}$ against $\sup _{t \in[0, T)} \operatorname{rank}\left(c_{t}\right) \neq r_{0}$ (which is much more useful). In the classical setting of i.i.d. or weakly dependent data, various estimation methods for the rank of an unknown covariance matrix (and related objects) have been proposed. We would like to mention Gaussian elimination method with complete pivoting of [3] and the test suggested in [13] among others. Unfortunately, these procedures can not be applied to our statistical problem as the probabilistic structure of the process $X$ is more complex and the rank is time-varying.

Our method is based upon a random perturbation of the original data and determinant expansions. The main idea can be described as follows: if we compute $\operatorname{det}\left(c_{t}+h e_{t}\right)$ for a positive definite $d \times d$ matrix $e_{t}$ independent of $c_{t}$ and $h \downarrow 0$, then, under appropriate conditions, its rate of decay to 0 depends on the unknown rank of $c_{t}$. Hence, the ratio $\operatorname{det}\left(c_{t}+2 h e_{t}\right) / \operatorname{det}\left(c_{t}+h e_{t}\right)$ asymptotically identifies the rank of $c_{t}$. Indeed, our main statistic is a partial sum of squared determinants of matrices built from $d$ consecutive increments of the process $X$ and the random perturbation is performed by a properly scaled Brownian motion $W^{\prime}$, which is independent of all ingredients of $X$. We remark that perturbation methods (and matrix expansions as well) find applications in various fields of mathematics; we refer, for instance, to [11] whose authors apply matrix perturbation methods to determine the number of components in a linear mixture model from high dimensional noisy samples. Furthermore, the methods of [2] also rely upon a generation of a new Brownian motion $W^{\prime}$ although in a completely different setting.

The paper is structured as follows. Section 2 is devoted to model assumptions, testing hypotheses and test statistics. We present the asymptotic theory for our estimators and apply it to maximal rank testing in Section 3. In Section 4, we develop a test for the null hypothesis of constant rank. All proofs are deferred to Section 6.

## 2. Model, assumptions and a random perturbation.

2.1. The setting and testing hypotheses. Our process of interest is a $d$-dimensional continuous Itô semimartingale $X$, given on some filtered probability space $\left(\Omega, \mathcal{F},\left(\mathcal{F}_{t}\right)_{t \geq 0}, \mathbb{P}\right)$. In vector form, and with $W$ denoting a $q$-dimensional Brownian motion, it can be written as

$$
\begin{equation*}
X_{t}=X_{0}+\int_{0}^{t} b_{s} d s+\int_{0}^{t} \sigma_{s} d W_{s} \tag{2.1}
\end{equation*}
$$

where $b_{t}$ is a $d$-dimensional drift process and $\sigma_{t}$ is a $\mathbb{R}^{d \times q}$-valued volatility process, assumed to be continuous in time [and indeed much more, see Assumption (H) below]. We set

$$
\begin{equation*}
c_{t}=\sigma_{t} \sigma_{t}^{*}, \quad r_{t}=\operatorname{rank}\left(c_{t}\right), \quad R_{t}=\sup _{s \in[0, t)} r_{s} \tag{2.2}
\end{equation*}
$$

We remark that the maximal rank $R_{T}$ is not bigger than the rank of the integrated volatility $\int_{0}^{T} c_{t} d t$, but may be strictly smaller. As already mentioned, it is suitable to use the smallest possible dimension for $W$, on the time interval $[0, T]$. This is the $\mathbb{P}$-essential supremum of $\omega \mapsto R_{T}(\omega)$, but, since a single path $t \mapsto X_{t}(\omega)$ is (partially) observed, the only available information is $R_{T}$ itself. So the problem really boils down to finding the behavior of the process $r_{t}$, and for this the choice of the dimension of $W$ in (2.1) is irrelevant.

The rank $r_{t}$ is the biggest integer $r \leq d$ such that the sum of the determinants of the matrices $\left(c_{t}^{i j}\right)_{i, j \in J}$, where $J$ runs through all subsets of $\{1, \ldots, d\}$ with $r$ points, is positive (with the convention that a $0 \times 0$ matrix has determinant 1 ) see, for example, [9], Lemma 3. Since $c_{t}$ is continuous, this implies that for any $r$ the random set $\left\{t: r_{t}(\omega)>r\right\}$ is open in $[0, T)$, so the mapping $t \mapsto r_{t}$ is lower semi-continuous. In particular, the set $\left\{t \in[0, T): r_{t}(\omega)=R_{T}(\omega)\right\}$ is a nonempty open subset. These properties also yield that the process $r_{t}$ is predictable and that the following subsets of $\Omega$, which later will be the "testing hypotheses," are $\mathcal{F}_{T}$-measurable:

$$
\begin{align*}
& \Omega_{T}^{r}=\left\{\omega: R_{T}(\omega)=r\right\}, \\
& \Omega_{T}^{=}=\left\{\omega: r_{t}(\omega)=R_{T}(\omega) \text { for all } t \in[0, T]\right\}, \\
& \Omega_{T}^{\neq}=\left\{\omega: t \mapsto r_{t}(\omega)\right. \text { has finitely many discontinuities and is, }  \tag{2.3}\\
& \qquad \quad \text { not Lebesgue-a.s. constant on }[0, T]\} .
\end{align*}
$$

Notice that we impose that $r_{T}=R_{T}$ in $\Omega_{\bar{T}}^{\overline{ }}$, whereas the lower semi-continuity only implies in general that $r_{T} \leq R_{T}$. Observe also that a priori $t \mapsto r_{t}$ may be Lebesgue-a.s. constant and still have discontinuities (even infinitely many) on $[0, T]$. So, $\Omega_{T}^{\overline{=}}$ and $\Omega_{T}^{\neq}$are disjoint but $\Omega_{T}^{\overline{=}} \cup \Omega_{T}^{\neq} \neq \Omega$ in general. The main aim of this paper is testing the null hypothesis $\Omega_{T}^{r}$ against $\Omega_{T}^{\neq r}=\bigcup_{r^{\prime} \neq r, 0 \leq r^{\prime} \leq d} \Omega_{T}^{r^{\prime}}$ (and related hypotheses) and testing the null hypothesis of $\Omega_{T}^{=}$against $\Omega_{T}^{\neq}$.
2.2. Matrix perturbation. In order to explain the main idea of our method, we need to introduce some notation. Recall that $d$ and $q$ are the dimensions of $X$ and $W$, respectively. Then $\mathcal{M}$ is the set of all $d \times d$ matrices, $\mathcal{M}_{r}$ for $r \in$ $\{0, \ldots, d\}$ is the set of all matrices in $\mathcal{M}$ with rank $r$, and $\mathcal{M}^{\prime}$ is the set of all $d \times q$ matrices. For any matrix $A$ we denote by $A_{i}$ the $i$ th column of $A$; for any vectors $x_{1}, \ldots, x_{d}$ in $\mathbb{R}^{d}$, we write $\operatorname{mat}\left(x_{1}, \ldots, x_{d}\right)$ for the matrix in $\mathcal{M}$ whose $i$ th column is the column vector $x_{i}$. For $r \in\{0, \ldots, d\}$ and $A, B \in \mathcal{M}$ we define

$$
\begin{equation*}
\mathcal{M}_{A, B}^{r}=\left\{G \in \mathcal{M}: G_{i}=A_{i} \text { or } G_{i}=B_{i} \text { with } \#\left\{i: G_{i}=A_{i}\right\}=r\right\} \tag{2.4}
\end{equation*}
$$

In other words, $\mathcal{M}_{A, B}^{r}$ is the set of all matrices $G \in \mathcal{M}$ with $r$ columns equal to those of $A$ (at the same places), and the remaining $d-r$ ones equal to those of $B$. Let us define

$$
\begin{equation*}
\gamma_{r}(A, B)=\sum_{G \in \mathcal{M}_{A, B}^{r}} \operatorname{det}(G) \tag{2.5}
\end{equation*}
$$

We demonstrate the main ideas for a deterministic problem first. Let $A \in \mathcal{M}$ be an unknown matrix with rank $r$. Assume that, although $A$ is unknown, we have a way of computing $\operatorname{det}(A+h B)$ for all $h>0$ and some given matrix $B \in \mathcal{M}_{d}$. We will see that the multi-linearity property of the determinant implies the following asymptotic expansion:

$$
\begin{equation*}
\operatorname{det}(A+h B)=h^{d-r} \gamma_{r}(A, B)+O\left(h^{d-r+1}\right) \tag{2.6}
\end{equation*}
$$

which is the core of our method. Thus, if $\gamma_{r}(A, B) \neq 0$, we have

$$
\begin{equation*}
\frac{\operatorname{det}(A+2 h B)}{\operatorname{det}(A+h B)} \rightarrow 2^{d-r} \quad \text { as } h \downarrow 0 \tag{2.7}
\end{equation*}
$$

and this convergence identifies the parameter $r$. However, it is impossible to choose a matrix $B \in \mathcal{M}$ which guarantees $\gamma_{r}(A, B) \neq 0$ for all $A \in \mathcal{M}_{r}$. To solve this problem, we can use a random perturbation. As we will show later, for any $A \in \mathcal{M}_{r}$ we have $\gamma_{r}(A, B) \neq 0$ a.s. when $B$ is the random matrix whose entries are independent standard normal. This idea will be the core of our testing procedure.
2.3. Assumptions and the test statistic. Before we proceed with the definition of the test statistic, we introduce the main assumptions. We need more structure than the mere equation (2.1), namely that the processes $b_{t}$ and $\sigma_{t}$, and also the volatility of $\sigma_{t}$, are continuous Itô semimartingales. In view of the previous discussion, it is no restriction to assume that all these are driven by the same $q$-dimensional Brownian motion, provided we take $q$ large enough. This leads us to put:

Assumption (H). The $d$-dimensional semimartingale $X$, defined on $(\Omega, \mathcal{F}$, $\left.\left(\mathcal{F}_{t}\right)_{t \geq 0}, \mathbb{P}\right)$, has the form

$$
\begin{align*}
X_{t} & =X_{0}+\int_{0}^{t} b_{s} d s+\int_{0}^{t} \sigma_{s} d W_{s} \\
\sigma_{t} & =\sigma_{0}+\int_{0}^{t} a_{s} d s+\int_{0}^{t} v_{s} d W_{s} \\
b_{t} & =b_{0}+\int_{0}^{t} a_{s}^{\prime} d s+\int_{0}^{t} v_{s}^{\prime} d W_{s}  \tag{2.8}\\
v_{t} & =v_{0}+\int_{0}^{t} a_{s}^{\prime \prime} d s+\int_{0}^{t} v_{s}^{\prime \prime} d W_{s}
\end{align*}
$$

where $W$ is a $q$-dimensional Brownian motion, and $b_{t}$ and $a_{t}^{\prime}$ are $\mathbb{R}^{d}$-valued, $\sigma_{t}, a_{t}$ and $v_{t}^{\prime}$ are $\mathbb{R}^{d \times q}$-valued, $v_{t}$ and $a_{t}^{\prime \prime}$ are $\mathbb{R}^{d \times q \times q}$-valued, and $v_{t}^{\prime \prime}$ is $\mathbb{R}^{d \times q \times q \times q}$-valued, all those processes being adapted. Finally, the processes $a_{t}$, $v_{t}^{\prime}, v_{t}^{\prime \prime}$ are càdlàg and the processes $a_{t}^{\prime}, a_{t}^{\prime \prime}$ are locally bounded.

At this stage it is not quite clear why the full force of Assumption (H) is required. In the standard limit theory for high frequency data of continuous Itô semimartingales, see, for example, [1, 7], only the first two representations of (2.8) are assumed. We will further explain Assumption (H) once we introduce the test statistic. When $b_{t}=g_{1}\left(X_{t}\right), \sigma_{t}=g_{2}\left(X_{t}\right)$ with $g_{1} \in C^{2}\left(\mathbb{R}^{d}\right)$ and $g_{2} \in C^{4}\left(\mathbb{R}^{d}\right)$, then Assumption (H) is automatically satisfied due to Itô's formula.

REMARK 2.1. Since $\sigma_{t}$ is not uniquely specified, whereas $c_{t}$ is, and since we really are interested in specific properties of $c_{t}$, it would be much nicer to replace the structural assumption on $\sigma_{t}$ [second equation in (2.8)] by a similar assumption on the process $c_{t}$ itself.

This is of course a trivial matter when $c_{t}$ is everywhere invertible: in this case $c_{t}$ is a continuous Itô semimartingale if and only if $\sigma_{t}$ is. But here we are precisely trying to describe the rank of the matrix $c_{t}$, so it is out of the question to assume that it is a priori invertible. Unfortunately, we were unable to replace the assumption on $\sigma$ by a similar (and de facto weaker) assumption on $c$.

Motivated by the matrix perturbation at (2.6), our tests will be based on statistics involving sums of (squared) determinants. The test function will be the nonnegative map $f$ on $\left(\mathbb{R}^{d}\right)^{d}$ defined as

$$
\begin{equation*}
f\left(x_{1}, \ldots, x_{d}\right)=\operatorname{det}\left(\operatorname{mat}\left(x_{1}, \ldots, x_{d}\right)\right)^{2} \tag{2.9}
\end{equation*}
$$

The authors of [9] used the following statistics:

$$
\begin{gather*}
\Delta_{n} \sum_{i=1}^{\left[t / \Delta_{n}\right]-d+1} f\left(\Delta_{i}^{n} X / \sqrt{\Delta_{n}}, \ldots, \Delta_{i+d-1}^{n} X / \sqrt{\Delta_{n}}\right)  \tag{2.10}\\
\Delta_{i}^{n} X=X_{i \Delta_{n}}-X_{(i-1) \Delta_{n}}
\end{gather*}
$$

to test for the full rank, thus allowing for efficient testing of the null hypothesis $\Omega_{T}^{d}$. On the sets $\Omega_{T}^{r}$ with $r<d$, however, it exhibits complex degeneracies and becomes difficult to study. In order to be able to analyze the asymptotic behavior of the preceding statistics, we introduce a random perturbation of the original process $X$ as motivated at the end of Section 2.2 (a somewhat similar idea in a different context was applied in [2]). More specifically, we choose a nonrandom invertible $d \times d$ matrix $\widetilde{\sigma}$ and generate a new process

$$
\begin{equation*}
X_{t}^{\prime}=\tilde{\sigma} W_{t}^{\prime} \tag{2.11}
\end{equation*}
$$

where $W^{\prime}$ is a $d$-dimensional Brownian motion independent of all processes in (2.8) [without loss of generality, for the mathematical treatment below we may assume that it is also defined on the space $\left.\left(\Omega, \mathcal{F},\left(\mathcal{F}_{t}\right)_{t \geq 0}, \mathbb{P}\right)\right]$. Following the ideas of Section 2.2, we add to $X$ this new process $X^{\prime}$, with a multiplicative factor going to 0 . As a matter of fact we introduce two such additions, and for $\kappa=1$ or 2 we set

$$
\begin{equation*}
Z_{t}^{n, \kappa}=X_{t}+\sqrt{\kappa \Delta_{n}} X_{t}^{\prime} \tag{2.12}
\end{equation*}
$$

Hence, with the notation of Section 2.2, we use $h=\sqrt{\Delta_{n}}$, which leads later to the optimal rate of convergence.

Another problem arises, namely in (2.10) successive summands partly use the same increments of $X$, and this causes problems for the central limit theorem. These problems can actually be overcome, at the expense of quite many additional technicalities, and with the advantage of a smaller asymptotic variance for our estimators below. However, in our case the crucial point is the choice of the tuning "parameter" $\tilde{\sigma}$ : this choice has an impact on the asymptotic variance as well, and since an "optimal" choice of $\widetilde{\sigma}$ seems out of reach, we will content ourselves with an arbitrary choice of $\widetilde{\sigma}$ and with a version of (2.10) with no overlapping of increments between the successive summands. This leads us to use the following two basic statistics:

$$
\begin{align*}
& S_{t}^{n, 1}=2 d \Delta_{n} \sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} f\left(\frac{Z_{(2 i d+1) \Delta_{n}}^{n, 1}-Z_{2 i d \Delta_{n}}^{n, 1}}{\sqrt{\Delta_{n}}}, \ldots,\right. \\
&\left.\frac{Z_{(2 i d+d) \Delta_{n}}^{n, 1}-Z_{(2 i d+d-1) \Delta_{n}}^{n, 1}}{\sqrt{\Delta_{n}}}\right), \\
& S_{t}^{n, 2}=2 d \Delta_{n} \sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} f\left(\frac{Z_{(2 i d+2) \Delta_{n}}^{n, 2}-Z_{(2 i d) \Delta_{n}}^{n, 2}}{\sqrt{2 \Delta_{n}}}, \ldots,\right.  \tag{2.13}\\
&\left.\frac{Z_{(2 i d+2 d) \Delta_{n}}^{n, 2}-Z_{(2 i d+2 d-2) \Delta_{n}}^{n, 2}}{\sqrt{2 \Delta_{n}}}\right) .
\end{align*}
$$

Notice that the statistics $S_{t}^{n, 1}$ and $S_{t}^{n, 2}$ are essentially the same (upon using $Z^{n, 1}$ and $Z^{n, 2}$, of course), except $S_{t}^{n, 2}$ is computed using the frequency $2 \Delta_{n}$.

At stage $n$ one observes the increments $\Delta_{i}^{n} X$ and simulates the increments $\Delta_{i}^{n} X^{\prime}$ for $i \leq\left[t / \Delta_{n}\right]$, so one "observes" all variables incurring in the definition of these two statistics.

Remark 2.2. Now, let us explain why Assumption (H) and the random perturbation in (2.12) are required. A direct stochastic expansion of the increments $\Delta_{i}^{n} Z^{n, 1}$ under Assumption (H) implies the decomposition

$$
\begin{align*}
& \operatorname{mat}\left(\Delta_{i}^{n} Z^{n} / \sqrt{\Delta_{n}}, \ldots, \Delta_{i+d-1}^{n} Z^{n} / \sqrt{\Delta_{n}}\right) \\
& \quad=\alpha_{i}^{n}+\sqrt{\Delta_{n}}\left(\beta_{i}^{n}(1)+\beta_{i}^{n}(2)\right)+O_{\mathbb{P}}\left(\Delta_{n}\right) \tag{2.14}
\end{align*}
$$

where the matrices $\alpha_{i}^{n}=\operatorname{mat}\left(\alpha_{i, 1}^{n}, \ldots, \alpha_{i, d}^{n}\right)$ and $\beta_{i}^{n}(k)=\operatorname{mat}\left(\beta_{i, 1}^{n}(k), \ldots, \beta_{i, d}^{n}(k)\right)$ in $\mathcal{M}$, for $k=1,2$, are given by

$$
\begin{align*}
\alpha_{i, j}^{n} & =\Delta_{n}^{-1 / 2} \sigma_{(i-1) \Delta_{n}} \Delta_{i+j-1}^{n} W \\
\beta_{i, j}^{n}(1) & =b_{(i-1) \Delta_{n}}+\Delta_{n}^{-1} v_{(i-1) \Delta_{n}} \int_{(i+j-1) \Delta_{n}}^{(i+j) \Delta_{n}}\left(W_{s}-W_{(i-1) \Delta_{n}}\right) d W_{s},  \tag{2.15}\\
\beta_{i, j}^{n}(2) & =\Delta_{n}^{-1 / 2} \widetilde{\sigma} \Delta_{i+j-1}^{n} W^{\prime} .
\end{align*}
$$

We remark that the matrices $\alpha_{i}^{n}, \beta_{i}^{n}(1), \beta_{i}^{n}(2)$ are $O_{\mathbb{P}}(1)$. In the case $r_{t} \leq d-1$ for all $t$, the first order term $\alpha_{i}^{n}$, which depends on the process $\sigma_{t}$, gives a degenerate limit when plugged in into the statistics (2.12) or (2.13). Hence the second order term $\sqrt{\Delta_{n}}\left(\beta_{i}^{n}(1)+\beta_{i}^{n}(2)\right)$, which involves the processes $b_{t}$ and $v_{t}$, becomes important. Indeed, we will see in Section 3 that it affects the limits. Furthermore, it is important to control the error of the above decomposition, and this is done by using the last two equations in (2.8).

The asymptotic expansion in (2.14) is a stochastic analogue of the perturbation presented in (2.6) (up to an error term) with $A=\alpha_{i}^{n}, B=\beta_{i}^{n}(1)+\beta_{i}^{n}$ (2) and $h=\sqrt{\Delta_{n}}$. Under Assumption (H) the term $\beta_{i}^{n}(1)$ already constitutes a random perturbation of the leading matrix $\alpha_{i}^{n}$. However, this perturbation does not guarantee that the quantity $\gamma_{r}\left(\alpha_{i}^{n}, \beta_{i}^{n}(1)\right)$ defined in (2.5) does not vanish when $\operatorname{rank}\left(\sigma_{(i-1) \Delta_{n}}\right)=r$ (which is essential for our method). To illustrate this problem, let us give a simple example. Let $d=3, q=1$ and define the processes

$$
d X_{t}^{j}=\sigma_{t}^{j} d W_{t}, \quad d \sigma_{t}^{j}=v_{t}^{j} d W_{t}, \quad j=1,2,3 .
$$

(So $W$ is a one-dimensional Brownian motion.) Then $\operatorname{rank}\left(\alpha_{i}^{n}\right)=1, \operatorname{rank}\left(\beta_{i}^{n}(1)\right)=$ 1 , and hence $\gamma_{1}\left(\alpha_{i}^{n}, \beta_{i}^{n}(1)\right)=0$. The presence of the new independent process $X^{\prime}$, and thus of the term $\beta_{i}^{n}(2)$, regularizes the problem. Indeed, we will show that $\gamma_{r}\left(\alpha_{i}^{n}, \beta_{i}^{n}(1)+\beta_{i}^{n}(2)\right)$ does not vanish whenever $\operatorname{rank}\left(\sigma_{(i-1) \Delta_{n}}\right)=r$. Finally, the perturbation rate $h=\sqrt{\Delta_{n}}$ in front of the process $X^{\prime}$ is chosen to achieve the best rate of convergence for the normalized versions of the statistics $S_{t}^{n, 1}, S_{t}^{n, 2}$.

Following the expansion (2.6), we know that the order of $\operatorname{det}\left(\alpha_{i}^{n}+\sqrt{\Delta_{n}}\left(\beta_{i}^{n}(1)+\right.\right.$ $\left.\left.\beta_{i}^{n}(2)\right)\right)^{2}$ is increasing in $r=\operatorname{rank}\left(\sigma_{(i-1) \Delta_{n}}\right)$. Consequently, as in (2.7), the ratio $S_{T}^{n, 2} / S_{T}^{n, 1}$ is expected to identify (asymptotically) the maximal rank $R_{T}$. The complete asymptotic theory is presented in the next section.

## 3. The asymptotic results and test for the maximal rank.

3.1. Notation. In order to present the main asymptotic results, we need to introduce some more notation. We define the function $F_{r}$ on $\left(\mathbb{R}^{2 d}\right)^{d}$ by

$$
\begin{align*}
& F_{r}\left(v_{1}, \ldots, v_{d}\right)=\gamma_{r}\left(\operatorname{mat}\left(x_{1}, \ldots, x_{d}\right), \operatorname{mat}\left(y_{1}, \ldots, y_{d}\right)\right)^{2} \\
& \qquad \text { if } v_{j}=\binom{x_{j}}{y_{j}} \in \mathbb{R}^{2 d} \tag{3.1}
\end{align*}
$$

where the quantity $\gamma_{r}$ is defined by (2.5). Next, let $\mathcal{U}=\mathcal{M}^{\prime} \times \mathcal{M} \times \mathbb{R}^{d q^{2}} \times \mathbb{R}^{d}$, whose points are $\underline{u}=(\alpha, \beta, \gamma, a)$, where $\alpha \in \mathcal{M}^{\prime}$ and $\beta \in \mathcal{M}$ and $\gamma \in \mathbb{R}^{d q^{2}}$ and $a \in \mathbb{R}^{d}$. Let us denote by $\bar{W}$ and $\bar{W}^{\prime}$ two independent Brownian motions with respective dimensions $q$ and $d$, defined on some space $\left(\bar{\Omega}, \overline{\mathcal{F}},\left(\overline{\mathcal{F}}_{t}\right), \overline{\mathbb{P}}\right)$. If $\underline{u} \in \mathcal{U}$ and $\kappa=1,2$ and $i \geq 1$, we associate the $2 d$-dimensional variables with the following components for $l=1, \ldots, d$ :

$$
\begin{align*}
\Psi(\underline{u}, \kappa)_{i}^{l}= & \frac{1}{\sqrt{\kappa}} \sum_{m=1}^{q} \alpha^{l m}\left(\bar{W}_{\kappa i}^{m}-\bar{W}_{\kappa(i-1)}^{m}\right), \\
\Psi(\underline{u}, \kappa)_{i}^{d+l}= & a^{l}+\frac{1}{\sqrt{\kappa}} \sum_{m=1}^{d} \beta^{l m}\left(\bar{W}_{\kappa i}^{m}-\bar{W}_{\kappa(i-1)}^{\prime m}\right)  \tag{3.2}\\
& +\frac{1}{\kappa} \sum_{m, k=1}^{q} \gamma^{l m k} \int_{\kappa(i-1)}^{\kappa i} \bar{W}_{s}^{k} d \bar{W}_{s}^{m} .
\end{align*}
$$

With the notation (3.1) we can then define the variables

$$
\begin{equation*}
\bar{F}_{r}(\underline{u}, \kappa)=F_{r}\left(\Psi(\underline{u}, \kappa)_{1}, \ldots, \Psi(\underline{u}, \kappa)_{d}\right) . \tag{3.3}
\end{equation*}
$$

The two sequences $(\Psi(\underline{u}, \kappa))_{i \geq 1}$ are not independent, but they have the same (global) law, for $\kappa=1,2$. Therefore if $\underline{u}=(\alpha, \beta, \gamma, a)$ we can set

$$
\begin{align*}
& \Gamma_{r}(\underline{u})=\overline{\mathbb{E}}\left(\bar{F}_{r}(\underline{u}, 1)\right)=\overline{\mathbb{E}}\left(\bar{F}_{r}(\underline{u}, 2)\right), \\
& \Gamma_{r}^{\prime}(\underline{u})=\overline{\mathbb{E}}\left(\bar{F}_{r}(\underline{u}, 1)^{2}\right)-\Gamma_{r}(\underline{u})^{2}=\overline{\mathbb{E}}\left(\bar{F}_{r}(\underline{u}, 2)^{2}\right)-\Gamma_{r}(\underline{u})^{2},  \tag{3.4}\\
& \Gamma_{r}^{\prime \prime}(\underline{u})=\overline{\mathbb{E}}\left(\bar{F}_{r}(\underline{u}, 1) \bar{F}_{r}(\underline{u}, 2)\right)-\Gamma_{r}(\underline{u})^{2} .
\end{align*}
$$

We then obtain the following crucial properties.

Lemma 3.1. Let $\underline{u}=(\alpha, \beta, \gamma, a) \in \mathcal{U}$ with $\beta \in \mathcal{M}_{d}$. Then if $r \in\{0,1, \ldots, d\}$,

$$
\begin{align*}
& \operatorname{rank}(\alpha)=r \quad \Longrightarrow \quad \Gamma_{r}(\underline{u})>0, \quad \Gamma_{r}^{\prime}(\underline{u})>\Gamma_{r}^{\prime \prime}(\underline{u}), \\
& \operatorname{rank}(\alpha)<r \quad \Longrightarrow \quad \Gamma_{r}(\underline{u})=\Gamma_{r}^{\prime}(\underline{u})=\Gamma_{r}^{\prime \prime}(\underline{u})=0 . \tag{3.5}
\end{align*}
$$

3.2. The limiting results. The key result is the asymptotic behavior of the processes $S^{n, j}$, which are defined by (2.13), as $n \rightarrow \infty$. These processes enjoy a law of large numbers and a central limit theorem, the centering being around one of the following processes, where $r$ is any (fixed) integer between 0 and $d$ :

$$
\begin{equation*}
S(r)_{t}=\int_{0}^{t} \Gamma_{r}\left(\sigma_{s}, \tilde{\sigma}, v_{s}, b_{s}\right) d s \tag{3.6}
\end{equation*}
$$

We will in fact have a CLT for the two-dimensional processes $U(r)^{n}$ with components

$$
\begin{equation*}
U(r)^{n, \kappa}=\frac{1}{\sqrt{\Delta_{n}}}\left(\frac{1}{\left(\kappa \Delta_{n}\right)^{d-r}} S^{n, \kappa}-S(r)\right) \tag{3.7}
\end{equation*}
$$

Of course, the centering process $S(r)$ depends on $r$, so one needs an additional assumption related with the particular value of $r$ which is chosen below (in contrast, the centering term is the same for all components).

THEOREM 3.2. Assume Assumption $(\mathrm{H})$, and also that $r_{t}(\omega) \leq r$ identically for some $r \in\{0, \ldots, d\}$. Then we have the stable (functional) convergence in law

$$
\begin{equation*}
U(r)^{n} \stackrel{\mathcal{L}-s}{\Longrightarrow} \mathcal{U}(r) \tag{3.8}
\end{equation*}
$$

where $\mathcal{U}(r)=\left(\mathcal{U}(r)^{\kappa}\right)_{\kappa=1,2}$ is defined on an extension $\left(\widetilde{\Omega}, \widetilde{\mathcal{F}},\left(\widetilde{\mathcal{F}}_{t}\right)_{t \geq 0}, \widetilde{\mathbb{P}}\right)$ of $\left(\Omega, \mathcal{F},\left(\mathcal{F}_{t}\right)_{t \geq 0}, \mathbb{P}\right)$ and is, conditionally on $\mathcal{F}$, a continuous centered Gaussian martingale with conditional covariance

$$
\begin{align*}
& \widetilde{\mathbb{E}}\left(\mathcal{U}(r)_{t}^{\kappa} \mathcal{U}(r)_{t}^{\kappa^{\prime}} \mid \mathcal{F}\right) \\
& \quad=V(r)_{t}^{\kappa \kappa^{\prime}}:= \begin{cases}2 d \int_{0}^{t} \Gamma_{r}^{\prime}\left(\sigma_{s}, \tilde{\sigma}, v_{s}, b_{s}\right) d s, & \text { if } \kappa=\kappa^{\prime}, \\
2 d \int_{0}^{t} \Gamma_{r}^{\prime \prime}\left(\sigma_{s}, \tilde{\sigma}, v_{s}, b_{s}\right) d s, & \text { if } \kappa \neq \kappa^{\prime} .\end{cases} \tag{3.9}
\end{align*}
$$

Note that in the above setting, if $r<r^{\prime} \leq d$, we also have $r_{t} \leq r^{\prime}$ and thus the results also hold with $r^{\prime}$ instead of $r$ everywhere. This does not bring a contradiction because, by (3.5), in this case the processes $S\left(r^{\prime}\right)$ and $U\left(r^{\prime}\right)$ are identically vanishing.

Now, these processes $S^{n, j}$ are only tools, and at the end we will be interested, for any $T>0$ fixed, in "estimators" for $R_{T}$, which are

$$
\begin{equation*}
\widehat{R}(n, T)=d-\frac{\log \left(S_{T}^{n, 2} / S_{T}^{n, 1}\right)}{\log 2} \tag{3.10}
\end{equation*}
$$

The quantity is a transformed analogue of the term on the left-hand side of (2.7). Since $R_{T}$ is by definition an integer between 0 and $d$, we can also use the following estimators:

$$
\widehat{R}^{\text {int }}(n, T)= \begin{cases}0, & \text { if } \widehat{R}(n, T)<\frac{1}{2}  \tag{3.11}\\ r, & \text { if } 1 \leq r \leq d-1, r-\frac{1}{2} \leq \widehat{R}(n, T)<r+\frac{1}{2} \\ d, & \text { if } \widehat{R}(n, T) \geq d-\frac{1}{2}\end{cases}
$$

The following corollary is a simple consequence of the previous theorem.
Corollary 3.3. Assume Assumption (H), and let $T>0$. Then the following stable convergence in law holds:

$$
\begin{equation*}
\frac{1}{\sqrt{\Delta_{n}}}(\widehat{R}(n, T)-r) \xrightarrow{\mathcal{L}-s} \mathcal{S}(T), \tag{3.12}
\end{equation*}
$$

where $\mathcal{S}(T)$ can be realized on the set $\Omega_{T}^{r}$, for any $r=0, \ldots, d$, as $\mathcal{S}(T)=$ $\frac{1}{\log 2}\left(\mathcal{U}(r)_{T}^{1}-\mathcal{U}(r)_{T}^{2}\right) / S(r)_{T}$ and is thus defined on an extension $\left(\widetilde{\Omega}, \widetilde{\mathcal{F}},\left(\widetilde{\mathcal{F}}_{t}\right)_{t \geq 0}, \widetilde{\mathbb{P}}\right)$ of $\left(\Omega, \mathcal{F},\left(\mathcal{F}_{t}\right)_{t \geq 0}, \mathbb{P}\right)$ and is, conditionally on $\mathcal{F}$, a centered Gaussian variable whose conditional variance is

$$
\widetilde{\mathbb{E}}\left((\mathcal{S}(T))^{2} \mid \mathcal{F}\right)=V(T)
$$

where $V(T)$ is a.s. positive and given by

$$
\begin{equation*}
V(T)=\frac{1}{(\log 2)^{2}} \frac{V(r)_{T}^{1,1}+V(r)_{T}^{2,2}-2 V(r)_{T}^{1,2}}{\left(S(r)_{T}\right)^{2}} \quad \text { on each set } \Omega_{T}^{r} \tag{3.13}
\end{equation*}
$$

In particular, for any $r=0, \ldots, d$, we have the following consistency result:

$$
\begin{equation*}
\mathbb{P}\left(\Omega_{T}^{r} \cap\left\{\widehat{R}^{\mathrm{int}}(n, T) \neq r\right\}\right) \rightarrow 0 \tag{3.14}
\end{equation*}
$$

In order to make this result feasible, we need consistent estimators for $V(T)$. For the denominator $S(r)_{T}^{2}$, we can of course take the square of $\Delta_{n}^{r-d} S_{T}^{n, 1}$. As for the numerator, we need estimators for $V(r)_{T}^{\kappa, \kappa^{\prime}}$. Up to normalization, natural ones are as follows:

$$
\begin{align*}
& V_{t}^{n, \kappa \kappa^{\prime}}=4 d^{2} \Delta_{n} \sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} f\left(\frac{Z_{(2 i d+\kappa) \Delta_{n}}^{n, \kappa}-Z_{2 i d \Delta_{n}}^{n, \kappa}}{\sqrt{\kappa \Delta_{n}}}, \ldots,\right. \\
&\left.\frac{Z_{(2 i d+\kappa d) \Delta_{n}}^{n, \kappa}-Z_{(2 i d+\kappa(d-1)) \Delta_{n}}^{n, \kappa}}{\sqrt{\kappa \Delta_{n}}}\right)  \tag{3.15}\\
& \times f\left(\frac{Z_{\left(2 i d+\kappa^{\prime}\right) \Delta_{n}}^{n, \kappa^{\prime}}-Z_{2 i d \Delta_{n}}^{n, \kappa^{\prime}}}{\sqrt{\kappa^{\prime} \Delta_{n}}}, \ldots,\right. \\
&\left.\frac{Z_{\left(2 i d+\kappa^{\prime} d\right) \Delta_{n}}^{n, \kappa^{\prime}}-Z_{\left(2 i d+\kappa^{\prime}(d-1)\right) \Delta_{n}}^{n, \kappa^{\prime}}}{\sqrt{\kappa^{\prime} \Delta_{n}}}\right),
\end{align*}
$$

where the function $f$ is given by (2.9).

## Proposition 3.4. Assume Assumption (H).

(a) If $r_{t}(\omega) \leq r$ identically for some $r \in\{0, \ldots, d\}$, we have for $\kappa, \kappa^{\prime}=1,2$ :

$$
\begin{align*}
\frac{1}{\left(\kappa \kappa^{\prime} \Delta_{n}^{2}\right)^{d-r}} V^{n, \kappa \kappa^{\prime}} \stackrel{\text { u.c.p. }}{\Longrightarrow} & 2 d \int_{0} \Theta_{s}^{r, \kappa, \kappa^{\prime}} d s  \tag{3.16}\\
\text { where } \Theta_{s}^{r, \kappa, \kappa^{\prime}}= & \left\{\begin{array}{c}
\Gamma_{r}^{\prime}\left(\sigma_{s}, \widetilde{\sigma}, v_{s}, b_{s}\right)+\Gamma_{r}\left(\sigma_{s}, \widetilde{\sigma}, v_{s}, b_{s}\right)^{2} \\
\text { if } \kappa=\kappa^{\prime} \\
\Gamma_{r}^{\prime \prime}\left(\sigma_{s}, \tilde{\sigma}, v_{s}, b_{s}\right)+\Gamma_{r}\left(\sigma_{s}, \tilde{\sigma}, v_{s}, b_{s}\right)^{2} \\
\text { if } \kappa \neq \kappa^{\prime},
\end{array}\right.
\end{align*}
$$

where the quantities $\Gamma_{r}, \Gamma_{r}^{\prime}, \Gamma_{r}^{\prime \prime}$ are defined in (3.4).
(b) We have

$$
\begin{equation*}
V(n, T):=\frac{V_{T}^{n, 11}+2^{2(\widehat{R}(n, T)-d)} V_{T}^{n, 22}-2^{1+\widehat{R}(n, T)-d} V_{T}^{n, 12}}{\left(S_{T}^{n, 1} \log 2\right)^{2}} \xrightarrow{\mathbb{P}} V(T) \tag{3.17}
\end{equation*}
$$

REMARK 3.5. In (3.17) one can replace $\widehat{R}(n, T)$ by $\widehat{R}^{\text {int }}(n, T)$. The numerator of the right-hand side of (3.13) is also $2\left(V(r)_{T}^{11}-V(r)_{T}^{12}\right)$. Therefore, we have

$$
V^{\prime}(n, T)=\frac{V_{T}^{n, 11}-2^{1+\widehat{R}(n, T)-d} V_{T}^{n, 12}}{\left(S_{T}^{n, 1}\right)^{2}} \xrightarrow{\mathbb{P}} V(T) \quad \text { on the set } \Omega_{T}^{r}
$$

as well. However, $V(n, T) \geq 0$ by construction (and it is even a.s. positive unless $r_{t}=0$ identically on $\left.[0, T]\right)$, a property not shared by $V^{\prime}(n, T)$.

Now, by the delta-method for stable convergence in law, the two previous results immediately yield the following corollary.

Corollary 3.6. Under Assumption (H) and for any $T>0$, we have

$$
\begin{equation*}
\frac{\widehat{R}(n, T)-R_{T}}{\sqrt{\Delta_{n} V(n, T)}} \xrightarrow{\mathcal{L}-s} \Phi, \tag{3.18}
\end{equation*}
$$

where $\Phi \sim \mathcal{N}(0,1)$ is defined on an extension $\left(\widetilde{\Omega}, \widetilde{\mathcal{F}},\left(\widetilde{\mathcal{F}}_{t}\right)_{t \geq 0}, \widetilde{\mathbb{P}}\right)$ of $(\Omega, \mathcal{F}$, $\left.\left(\mathcal{F}_{t}\right)_{t \geq 0}, \mathbb{P}\right)$ and is independent of $\mathcal{F}$.
3.3. Tests for the maximal rank. So far, $\widehat{R}(n, T)$ are estimators for the maximal rank $R_{T}$, which equals $r$ on the set $\Omega_{T}^{r}$, and even feasible estimators if we use Corollary 3.6. In particular, this corollary seems to allow us to easily construct confidence intervals for $R_{T}$. However, confidence interval do not make much statistical sense, since $R_{T}$ is integer-valued.

One might also use, perhaps more to the point, $\widehat{R}^{\text {int }}(n, T)$, but apart from the consistency result (3.14) it is not easy to evaluate the error probability that $\widehat{R}^{\text {int }}(n, T) \neq R_{T}$.

So, it seems more appropriate here to do testing: we can test the null hypothesis that the path lies in $\Omega_{T}^{r}$ for some $r$, against the alternative that it is in $\Omega_{T}^{r^{\prime}}$ for another specific $r^{\prime} \neq r$, or for all $r^{\prime}>r$ or all $r^{\prime}<r$, or all $r^{\prime} \neq r$. We may also use composite null hypotheses, such as being in $\Omega_{T}^{r}$ for some $r$ smaller, or bigger, than a given value $r_{0}$.

We start with the problem of testing the null hypothesis $\Omega_{T}^{r}$, against the alternative $\Omega_{T}^{\neq r}=\bigcup_{r^{\prime} \neq r, 0 \leq r^{\prime} \leq d} \Omega_{T}^{r^{\prime}}$. For any $\alpha \in(0,1)$, and with $z_{\alpha}$ being the symmetric $\alpha$-quantile of $\mathcal{N}(0,1)$ defined by $\mathbb{P}\left(|\Phi|>z_{\alpha}\right)=\alpha$ when $\Phi \sim \mathcal{N}(0,1)$, we take the critical (rejection) region

$$
\begin{equation*}
\mathcal{C}(\alpha)_{T}^{n,=r}=\left\{\omega:|\widehat{R}(n, T)-r|>z_{\alpha} \sqrt{\Delta_{n} V(n, T)}\right\} . \tag{3.19}
\end{equation*}
$$

Proposition 3.7. Under Assumption (H), the tests (3.19) have the asymptotic level $\alpha$ for testing the null $\Omega_{T}^{r}$, in the sense that

$$
\begin{equation*}
A \subset \Omega_{T}^{r}, \quad \mathbb{P}(A)>0 \Rightarrow \mathbb{P}\left(\mathcal{C}(\alpha)_{T}^{n,=r} \mid A\right) \rightarrow \alpha \tag{3.20}
\end{equation*}
$$

[above, $\mathbb{P}(\cdot \mid A)$ is the usual conditional probability]. They are also consistent for the alternative $\Omega_{T}^{\neq r}$, in the sense that

$$
\begin{equation*}
\mathbb{P}\left(\mathcal{C}(\alpha)_{T}^{n,=r} \cap \Omega_{T}^{\neq r}\right) \rightarrow \mathbb{P}\left(\Omega_{T}^{\neq r}\right) \tag{3.21}
\end{equation*}
$$

One constructs one-sided tests in the same way. For example, if we want to test the null hypothesis $\Omega_{T}^{\leq r}=\bigcup_{r^{\prime} \leq r} \Omega_{T}^{r^{\prime}}$ against the alternative $\Omega_{T}^{>r}=\bigcup_{r^{\prime}>r} \Omega_{T}^{r^{\prime}}$, and if $z_{\alpha}^{\prime}$ is the one-sided $\alpha$-quantile defined by $\mathbb{P}\left(\Phi>z_{\alpha}^{\prime}\right)=\alpha$, we take the critical region

$$
\begin{equation*}
\mathcal{C}(\alpha)_{T}^{n, \leq r}=\left\{\omega: \widehat{R}(n, T)>r+z_{\alpha}^{\prime} \sqrt{\Delta_{n} V(n, T)}\right\} . \tag{3.22}
\end{equation*}
$$

Exactly as above, one obtains the following proposition.
Proposition 3.8. Under Assumption (H), the tests (3.22) have the asymptotic level at most $\alpha$ for testing the null $\Omega_{T}^{\leq r}$, and indeed satisfy

$$
\begin{equation*}
A \subset \Omega_{T}^{\leq r}, \quad \mathbb{P}(A)>0 \Rightarrow \mathbb{P}\left(\mathcal{C}(\alpha)_{T}^{n, \leq r} \mid A\right) \rightarrow \alpha \mathbb{P}\left(\Omega_{T}^{\leq r} \mid A\right) \leq \alpha \tag{3.23}
\end{equation*}
$$

and are consistent for the alternative $\Omega_{T}^{>r}$.
The tests for the null $\Omega_{T}^{\geq r}$ against $\Omega_{T}^{<r}$ are obtained analogously.
REMARK 3.9. Apart from the direct rank testing, let us link our testing procedure with some other statistical problems:
(a) In financial mathematics the full rank of the volatility matrix $c_{t}$ is required for market completeness. Thus, testing $\Omega_{T}^{\leq d}$ versus $\Omega_{T}^{=d}$ is directly connected to completeness of the market.
(b) In $[5,6]$ parametric estimation methods for the so called integrated diffusions have been developed. An integrated diffusion is a process that satisfies the first and the third equations of Assumption (H) with $\sigma=0$, that is,

$$
d X_{t}=b_{t} d t
$$

where $b_{t}$ is a continuous Itô semimartingale. We refer to [5] for various applications of these models in natural sciences. Given high frequency observations of $X$, testing the null hypothesis of integrated diffusion versus the alternative of a diffusion with a present volatility part $\sigma$ is equivalent to testing $\Omega_{T}^{0}$ versus $\Omega_{T}^{\neq 0}$.
(c) Another potential application of our method is a test for "perfect correlation" between the process $X$ and the unobserved volatility $\sigma$. The problem can be formulated as follows: Let $X$ and $\sigma$ be two one-dimensional continuous Itô semimartingales of the form

$$
d X_{t}=b_{t} d t+\sigma_{t} d W_{t}, \quad d \sigma_{t}=a_{t} d t+v_{t} d B_{t}
$$

where $W$ and $B$ are one-dimensional Brownian motions with the bracket process $[W, B]_{t}=\rho t,|\rho| \leq 1$. For financial applications testing, the hypothesis $|\rho|=1$ versus $|\rho|<1$ is of certain interest. Note that $|\rho|=1$ appears in the SDE case, that is, when $\sigma_{t}=g\left(X_{t}\right)$ with $g \in C^{2}(\mathbb{R})$. We refer to testing local volatility hypothesis in [12] for a more detailed discussion (see also [14] for related statistical problems). The aforementioned problem is equivalent to testing $\Omega_{T}^{1}$ versus $\Omega_{T}^{>1}=\Omega_{T}^{2}$ for the two-dimensional process $(X, \sigma)$. Since the process $\sigma$ is unobserved, it has to be locally estimated from the high frequency observations of $X$ first (see, e.g., [12] for more details).
4. A test for a constant rank. This section is devoted to a seemingly different topic, namely whether the a priori time-dependent rank is constant or not. Our test statistics will be based on a distance measure between the rank process $r_{t}$ and the maximal rank $R_{T}$, which vanishes if and only if the rank is constant almost surely. For the formal testing procedure, we will need some limiting results for the "spot estimators" of the rank. By this, we mean estimators for $r_{t}$, for any given $t$, at least under the assumption that $r_{s}$ is equal to $r_{t}$ for all $s$ in some right or left neighborhood of $t$.

To describe these spot estimators, we pick a sequence $k_{n} \geq 1$ of integers going to infinity, and such that $k_{n} \Delta_{n} \rightarrow 0$ (as for spot volatility estimators), and precise specifications for $k_{n}$ will be given later, although we always assume $k_{n} \geq 4 d$. For any integer $i \geq 1$, we set

$$
\begin{equation*}
\widehat{R}_{i}^{n}=d-\frac{\log \widehat{S}_{i}^{n}}{\log 2}, \quad \widehat{S}_{i}^{n}=\frac{S_{2 d(i+1) k_{n} \Delta_{n}}^{n, 2}-S_{2 d i k_{n} \Delta_{n}}^{n, 2}}{S_{2 d(i+1) k_{n} \Delta_{n}}^{n, 1}-S_{2 d i k_{n} \Delta_{n}}^{n, 1}} \tag{4.1}
\end{equation*}
$$

Then $\widehat{R}_{i}^{n}$, more or less, plays the role of an estimator of the maximum of $r_{t}$ over an interval of length $2 d k_{n} \Delta_{n}$ around the time $2 i d \Delta_{n}$, and we set for any $p>0$ :

$$
\begin{align*}
c A(p)_{t}^{n} & =2 d k_{n} \Delta_{n} \sum_{i=0}^{\left[t / 2 d k_{n} \Delta_{n}\right]-2}\left\{\left|\widehat{R}_{i k_{n}}^{n}\right|^{p} \wedge(d+1)^{p}\right\}, \\
B(n, p, T) & =A(p)_{T}^{n}-a(n, T)(\widehat{R}(n, T))^{p},  \tag{4.2}\\
a(n, T) & =2 d k_{n} \Delta_{n}\left(\left[T / 2 d k_{n} \Delta_{n}\right]-1\right) .
\end{align*}
$$

The asymptotic results for the quantity $B(n, p, T)$ are as follows.
THEOREM 4.1. Assume Assumption (H), and let $T>0, p>0$ and $k_{n}$ be such that $k_{n} \Delta_{n}^{3 / 4} \rightarrow \infty$ and $k_{n} \Delta_{n} \rightarrow 0$.
(a) If $t \mapsto r_{t}(\omega)$ is continuous except at finitely many points on $[0, T]$, hence piecewise constant, we have

$$
\begin{equation*}
B(n, p, T) \xrightarrow{\mathbb{P}} \int_{0}^{T}\left(r_{s}\right)^{p} d s-T\left(R_{T}\right)^{p} \tag{4.3}
\end{equation*}
$$

(b) We have the stable convergence in law:

$$
\begin{equation*}
\frac{1}{\sqrt{\Delta_{n}}} B(n, p, T) \xrightarrow{\mathcal{L}-s} \mathcal{B}(p, T) \quad \text { in restriction to the set } \Omega_{T}^{\overline{=}} \cap\left\{R_{T} \geq 1\right\} \tag{4.4}
\end{equation*}
$$

where $\mathcal{B}(p, T)$ is defined on an extension $\left(\widetilde{\Omega}, \widetilde{\mathcal{F}},\left(\widetilde{\mathcal{F}}_{t}\right)_{t \geq 0}, \widetilde{\mathbb{P}}\right)$ of $\left(\Omega, \mathcal{F},\left(\mathcal{F}_{t}\right)_{t \geq 0}, \mathbb{P}\right)$ and is, conditionally on $\mathcal{F}$, a centered Gaussian variable with conditional variance $\bar{V}(p, T)=\widetilde{\mathbb{E}}\left(\mathcal{B}(p, T)^{2} \mid \mathcal{F}\right)$ given on each set $\Omega_{T}^{r}$ by

$$
\begin{align*}
\bar{V}(p, T)=\left(\frac{p r^{p-1}}{\log 2}\right)^{2} \int_{0}^{T} & \left(\frac{1}{\Gamma_{r}\left(\sigma_{s}, \tilde{\sigma}, v_{s}, b_{s}\right)}-\frac{T}{S(r)_{T}}\right)^{2}  \tag{4.5}\\
& \times\left(d V(r)_{s}^{11}+d V(r)_{s}^{22}-2 d V(r)_{s}^{12}\right)
\end{align*}
$$

with $V(r)^{\kappa \kappa^{\prime}}$ being defined at (3.9).
Notice that the right-hand side of (4.3) is 0 on the set $\Omega_{T}^{\overline{=}}$, and strictly negative on $\Omega_{T}^{\neq}$.

REMARK 4.2. The reader will notice that in the definition of $A(p)_{t}^{n}$ the summands are $\left|\widehat{R}_{i}^{n}\right|^{p} \wedge(d+1)^{d}$, instead of the more natural $\left|\widehat{R}_{i}^{n}\right|^{p}$. We could take this more natural form for (b) above, but it is useful (and innocuous from a practical viewpoint) to "bound" the summands, in order to obtain (a). We could bound them by $d^{p}$ instead of $(d+1)^{p}$ and still have (4.3), but then (4.4) would then fail in case $r=d$ is the maximal rank: we would obtain a CLT with a non-Gaussian and noncentered limit.

REMARK 4.3. In the setting of (b) above, we will in fact prove a joint convergence for the variables $A(p)_{T}^{n}-a(n, T) r^{p}$ and $\widehat{R}(n, T)-r$, both normalized by $1 / \sqrt{\Delta_{n}}$ [the second one being as in (3.12)], and from which (4.4) follows. Such a joint CLT even holds under the assumptions of (a), with a complicated limit, but this refinement is not useful for us in this paper.

REMARK 4.4. One can also prove a joint convergence for the variables $A(p)_{T}^{n}-a(n, T) r^{p}$ with different values of $p$, and still normalized by $1 / \sqrt{\Delta_{n}}$. However, when $p>p^{\prime}>0$ it turns out that the difference $\frac{1}{\sqrt{\Delta_{n}}}\left(A(p)_{T}^{n}-\right.$ $\left.a(n, T)^{1-p^{\prime} / p}\left(A\left(p^{\prime}\right)_{T}^{n}\right)^{p / p^{\prime}}\right)$ converges to 0 , and no known normalization gives a proper CLT.

As before, we need consistent estimators for the conditional variance $\bar{V}(p, T)$. Such estimators are constructed in a way analogous to (3.15). That is, we set with $k_{n}$ as above:

$$
\begin{align*}
& \bar{V}_{t}^{n, \kappa \kappa^{\prime}}=4 d^{2} \Delta_{n}^{1+2 d-2 \widehat{R}(n, T)} \\
& \times \sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-k_{n}-1}( \left.\frac{2 d k_{n} \Delta_{n}}{S_{2 d\left(i+k_{n}\right) \Delta_{n}}^{n, 1}-S_{2 i d \Delta_{n}}^{n, 1}}-\frac{T}{S_{T}^{n, 1}}\right)^{2} \\
& \times f\left(\frac{Z_{(2 i d+\kappa) \Delta_{n}}^{n, \kappa}-Z_{2 i d \Delta_{n}}^{n, \kappa}}{\sqrt{\kappa \Delta_{n}}}, \ldots,\right. \\
&\left.\frac{Z_{(2 i d+\kappa d) \Delta_{n}}^{n, \kappa}-Z_{(2 i d+\kappa(d-1)) \Delta_{n}}^{n, \kappa}}{\sqrt{\kappa \Delta_{n}}}\right)  \tag{4.6}\\
& \times f\left(\frac{Z_{\left(2 i d+\kappa^{\prime}\right) \Delta_{n}}^{n, \kappa^{\prime}}-Z_{2 i d \Delta_{n}}^{n, \kappa^{\prime}}}{\sqrt{\kappa^{\prime} \Delta_{n}}}, \ldots,\right. \\
&\left.\frac{Z_{\left(2 i d+\kappa^{\prime} d\right) \Delta_{n}}^{n, \kappa^{\prime}}-Z_{\left(2 i d+\kappa^{\prime}(d-1)\right) \Delta_{n}}^{n, \kappa^{\prime}}}{\sqrt{\kappa^{\prime} \Delta_{n}}}\right) .
\end{align*}
$$

THEOREM 4.5. Assume Assumption (H), and let $T>0, p>0$ and $k_{n}$ be such that $k_{n} \Delta_{n}^{3 / 4} \rightarrow \infty$ and $k_{n} \Delta_{n} \rightarrow 0$. Then we have

$$
\begin{align*}
& \bar{V}(n, p, T):=\left(\frac{p \widehat{R}(n, T)^{p-1}}{\log 2}\right)^{2} \\
& \times\left(\bar{V}_{T}^{n, 11}+2^{2(\widehat{R}(n, T)-d)} \bar{V}_{T}^{n, 22}-2^{1+\widehat{R}(n, T)-d} \bar{V}_{T}^{n, 12}\right)  \tag{4.7}\\
& \xrightarrow{\mathbb{P}} \bar{V}(p, T) \quad \text { on the set } \Omega_{\bar{T}}^{\overline{\bar{C}}} .
\end{align*}
$$

## Moreover, the variables

$$
\begin{equation*}
Z(n, p, T)=\frac{B(n, p, T)}{\left.\sqrt{\Delta_{n}\left(\bar{V}(n, p, T) \wedge\left(1 / \sqrt{\Delta_{n}}\right)\right.}\right)} \tag{4.8}
\end{equation*}
$$

have the following asymptotic behavior, where $\Phi \sim \mathcal{N}(0,1)$ is as in Corollary 3.6:

$$
\begin{array}{ll}
Z(n, p, T) \xrightarrow{\mathcal{L}-s} \Phi & \text { in restriction to the set } \Omega_{T}^{\overline{=}} \cap\left\{R_{T} \geq 1\right\}, \\
Z(n, p, T) \xrightarrow{\mathbb{P}}-\infty & \text { in restriction to the set } \Omega_{T}^{\neq} . \tag{4.9}
\end{array}
$$

Having all instruments at hand we proceed with testing. What is easily available is a family of tests for the null $\Omega_{\bar{T}}^{\overline{=}}$, whereas the alternative is restricted to $\Omega_{T}^{\neq}$. One does not know how to test the null $\Omega_{T}^{\neq}$.

For this purpose, we use the statistic $B(n, p, T)$. In fact, (4.9) gives us the behavior of this statistic on $\Omega_{T}^{=} \cap\left\{R_{T} \geq 1\right\}$, and this is the null which is tested below. Now, $\Omega_{T}^{=}$is the union of $\Omega_{T}^{=} \cap\left\{R_{T} \geq 1\right\}$ and $\Omega_{T}^{0}$, so if we are interested in testing the whole $\Omega_{T}^{\overline{=}}$ one can do a double test, using what precedes and Proposition 3.8 with $r=0$.

We propose to use the following critical region, where $p>0$ is chosen arbitrarily and $z_{\alpha}^{\prime}$ is again the one-sided $\alpha$-quantile of $\mathcal{N}(0,1)$ :

$$
\begin{equation*}
\mathcal{C}(\alpha)_{T}^{n, \equiv}=\left\{\omega: B(n, p, T)<-z_{\alpha}^{\prime} \sqrt{\Delta_{n}\left(\bar{V}(n, p, T) \wedge\left(1 / \sqrt{\Delta_{n}}\right)\right.}\right\} . \tag{4.10}
\end{equation*}
$$

Exactly as in the previous section, we obtain the following result.
Proposition 4.6. Under Assumption (H), the tests (4.10) have the asymptotic level $\alpha$ for testing the null $\Omega_{\bar{T}}^{\overline{=}} \cap\left\{R_{T} \geq 1\right\}$, in the sense of (3.20), and are consistent for the alternative $\Omega_{T}^{\neq}$.
5. Monte Carlo evidence. We have tested the performances of the test described in Section 3.3 through some Monte Carlo experiments. These are performed in the 2-dimensional case $d=2$, and first we conducted preliminary experiments to determine a reasonable value of $\tilde{\sigma}$ in (2.11). To this aim, we simulated the processes over the time interval $[0,1]$ and $\Delta_{n}=1 / 25,000$ (corresponding approximately to 1 second for daily observations). In all cases, $X$ has no drift and constant volatility, and $\widetilde{\sigma}=\lambda I_{d}$ and we considered two cases:

1. $X^{1}=W^{1}$ and $X^{2}=0$, so the rank is identically $r=1$;
2. $X^{1}=W^{1}$ and $X^{2}=W^{2}$, so the rank is identically $r=2$;
(the case of rank 0 , i.e., $X=0$ here, is left out since in this case the value of the constant $\lambda$ is irrelevant, by scaling). We measure the performance of the procedure by reporting, as a function of $\lambda$, the empirical standard deviation of $\widehat{R}(n, 1)-r$, over 5000 simulations:

| $\boldsymbol{\lambda}$ | $\mathbf{0 . 0 1 2}$ | $\mathbf{0 . 0 3 7}$ | $\mathbf{0 . 1 1}$ | $\mathbf{0 . 3 3}$ | $\mathbf{1}$ | $\mathbf{3}$ | $\mathbf{9}$ | $\mathbf{2 7}$ | $\mathbf{8 1}$ | $\mathbf{2 4 3}$ |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| stdv, $r=1$ | 0.053 | 0.052 | 0.051 | 0.052 | 0.052 | 0.052 | 0.052 | 0.064 | 0.27 | 0.76 |
| stdv, $r=2$ | 0.052 | 0.052 | 0.052 | 0.051 | 0.052 | 0.051 | 0.052 | 0.095 | 0.53 | 1.52 |

The results are not sensitive to the value of $\lambda$, as long as $\lambda$ is smaller than 10 : so in practice it seems advisable to take for $\tilde{\sigma}$ the identity matrix times the average value of the volatility of $X$ (i.e., when the two components of $X$ have the same order of variability).

Next, we conduct the Monte Carlo, again over 5000 simulations, in the case of a time-varying but nonstochastic volatility, and with $\lambda=2$. We think that a stochastic volatility does not really affect the problem at hand, and taking always the same deterministic volatility allows for better (or at least, easier) comparison between the different simulations. The time interval is still $[0,1]$ with $\Delta_{n}=1 / 25,000$, and we consider the function

$$
f(t)=1+(2 t-1)^{2}
$$

(which decreases from 2 to 1 and increases back to 2 when $t$ goes from 0 to 1). Then we again consider two cases:

- Case 1: the matrix $\sigma_{t}$ is diagonal with entries

$$
\sigma_{t}^{11}=f(t), \quad \sigma_{t}^{22}=f(t) 1_{\{t \leq S\}}
$$

for some $S \in[0,1]$, so $r_{t}=2$ if $t<S$ and $r_{t}=1$ if $t>S$, and thus $R_{1}=2$ if $S>0$ and $R_{1}=1$ if $S=0$. We also added a rather large drift with components $b_{t}^{1}=b_{t}^{2}=3$.

- Case 2: the volatility matrix $\sigma_{t}$ and the drift vector are

$$
\begin{aligned}
\sigma_{t} & =f(t)\left(\begin{array}{cc}
\cos (t A \pi / 2) & \cos (t A \pi / 2) \\
\sin (t A \pi / 2) & \sin (t A \pi / 2)
\end{array}\right) \\
b_{t} & =b^{\prime}\binom{1+\sin (t A \pi / 2)}{1+\cos (t A \pi / 2)}
\end{aligned}
$$

for some constant $b^{\prime}$, so $r_{t}=1$ identically, although the structure of the drift tends to "increase" the apparent rank when $A>0$. Unless $A=0$ (in which case the continuous martingale part of $X^{2}$ is identically vanishing), this implies that the nonnull eigenvector of $\sigma_{t}$ (and of $c_{t}$ as well) rotates with the period $4 / A$, so when $A$ increases the rotation becomes faster.

Below, in each table we report five quantities: the second and fourth empirical moments of the studentized variables $\frac{\widehat{R}(n, 1)-r}{\sqrt{\Delta_{n} V(n, 1)}}$ (to check for normality); next, the proportion of rejection for the three possible null hypotheses $R_{1}=0, R_{1}=1$ and $R_{1}=2$, for a nominal level $\alpha=0.05$.

In Case 1 we have the following results:

| $\boldsymbol{S}$ | 2d moment | 4th moment | Test $\boldsymbol{R}_{\mathbf{1}}=\mathbf{0}$ | Test $\boldsymbol{R}_{\mathbf{1}}=\mathbf{1}$ | Test $\boldsymbol{R}_{\mathbf{1}}=\mathbf{2}$ |
| :--- | :---: | :---: | :---: | :---: | :---: |
| 0.0 | 1.04 | 3.34 | 1.00 | 0.053 | 1.000 |
| 0.1 | 1.02 | 3.29 | 1.00 | 1.000 | 0.050 |
| 0.2 | 1.05 | 3.29 | 1.00 | 1.000 | 0.051 |
| 0.5 | 1.00 | 3.06 | 1.00 | 1.000 | 0.049 |
| 0.8 | 0.98 | 2.87 | 1.00 | 1.000 | 0.050 |

(recall that when $S=0$ we have $R_{1}=1$, and otherwise $R_{1}=2$ ). The columns giving moments point toward a fairly good normal approximation, the actual level of the tests for the correct value of $R_{1}$ is rather close to the nominal value 0.05 , and the power is always 1 : this is comforted by the fact that when we use the estimator $\widehat{R}^{\text {int }}(n, 1)$, the number of wrong decisions (when the estimator is not the correct value) is overall 4 , for $10 \times 5000$ simulations.

In Case 2 we have performed the simulations for various values of the constants $b^{\prime}$ and $A$. The results are fairly stable when these constants vary, so we only report three cases for $A$ and three cases for $b^{\prime}$ as well.

| $\boldsymbol{b}^{\prime}$ | $\boldsymbol{A}$ | $\boldsymbol{2 d}$ moment | 4th moment | Test $\boldsymbol{R}_{\mathbf{1}}=\mathbf{0}$ | Test $\boldsymbol{R}_{\mathbf{1}}=\mathbf{1}$ | Test $\boldsymbol{R}_{\mathbf{1}}=\mathbf{2}$ |
| ---: | ---: | :---: | :---: | :---: | :---: | :---: |
| 0 | 0 | 1.01 | 2.99 | 1.00 | 0.050 | 1.00 |
| 0 | 5 | 1.00 | 2.94 | 1.00 | 0.049 | 1.00 |
| 0 | 10 | 1.01 | 3.36 | 1.00 | 0.062 | 1.00 |
| 3 | 0 | 1.02 | 3.26 | 1.00 | 0.050 | 1.00 |
| 3 | 5 | 1.02 | 3.20 | 1.00 | 0.052 | 1.00 |
| 3 | 10 | 1.01 | 3.03 | 1.00 | 0.049 | 1.00 |
| 12 | 0 | 0.98 | 2.88 | 1.00 | 0.049 | 1.00 |
| 12 | 3 | 1.02 | 3.13 | 1.00 | 0.052 | 1.00 |
| 12 | 10 | 1.01 | 2.99 | 1.00 | 0.050 | 1.00 |

Comments analogous to those for Case 1 can be made, and here the number of wrong decisions turns out to be 0 for the whole set of simulations.
6. Proofs. Before we start presenting the formal proofs, let us give the road map. Section 6.1 demonstrates some technical results on expansions of determinants. They are applied in Section 6.2 to prove Lemma 3.1. This lemma implies that the process $S(r)_{t}$ defined at (3.6) is strictly positive on the set $\Omega_{T}^{r}$, which is crucial for our method.

The first main result of our paper is Theorem 3.2 whose proof is rather involved. First, we will show that the standard localization procedure (see, e.g., Section 3 in [1]) implies that all processes in Assumption (H) may be assumed to
be bounded without loss of generality. This first step considerably simplifies the stochastic treatment of various quantities. A second crucial step is the stochastic expansion explained in Remark 2.2: we have (2.14) and (2.15). Section 6.3 deals with the formal justification of this expansion, for which we will use slightly different notation.

It turns out that the stochastic order of the error term related to the decomposition (2.14), namely $O_{\mathbb{P}}\left(\Delta_{n}\right)$, is not sufficient to show its asymptotic negligibility. However, we will prove that the error terms are martingale differences, so they will not affect the stable central limit theorem at (3.8). A similar treatment will be required for the error term connected with the stochastic version of the expansion (2.6).

The proof of Proposition 3.4 (consistent estimation of the asymptotic conditional covariance matrix) is somewhat easier. Corollary 3.3 follows essentially from Theorem 3.2 by the delta method for stable convergence. The proofs of these results are collected in Section 6.4. In particular, we apply a stable central limit theorem for semimartingales (see, e.g., [8], Theorem IX.7.28) to prove Theorem 3.2.

The proofs of Theorems 4.1 and 4.5, which are presented in Section 6.5, are a bit more involved than that of Theorem 3.2, although the main techniques are similar. The additional difficulty comes from the fact that we need to use the stable convergence of Theorem 3.2, but for processes evaluated at random times. Corollary 3.6 and Propositions 3.7, 3.8 and 4.6 are straightforward consequences of the previous results.
6.1. Expansion of determinants. We first prove some general and easy facts about determinants. Below $\|A\|$ denotes the Euclidean norm of a matrix $A \in \mathcal{M}$.

For $m \geq 1$, we call $\mathcal{P}_{m}$ the set of all multi-integers $\mathbf{p}=\left(p_{1}, \ldots, p_{m}\right)$ with $p_{1}+$ $\cdots+p_{m}=d$, and $\mathcal{I}_{\mathbf{p}}$ is the set of all partitions $\mathbf{I}=\left(I_{1}, \ldots, I_{m}\right)$ of $\{1, \ldots, d\}$ such that $I_{j}$ contains exactly $p_{j}$ points (so $I_{j}=\varnothing$ if $p_{j}=0$ ). If $\mathbf{p} \in \mathcal{P}_{m}$ and $\mathbf{I} \in \mathcal{I}_{\mathbf{p}}$ and $A_{1}, \ldots, A_{m} \in \mathcal{M}$, we write $G_{A_{1}, \ldots, A_{m}}^{\mathbf{I}}$ for the matrix whose $i$ th column is the $i$ th column of $A_{j}$ when $i \in I_{j}$. Letting $A, B, C \in \mathcal{M}$, we can rewrite (2.5) as

$$
\begin{equation*}
\gamma_{r}(A, B)=\sum_{\mathbf{I} \in \mathcal{I}_{(r, d-r)}} \operatorname{det}\left(G_{A, B}^{\mathbf{I}}\right) \tag{6.1}
\end{equation*}
$$

and we set

$$
\begin{equation*}
\gamma_{r}^{\prime}(A, B, C)=\sum_{\mathbf{I} \in \mathcal{I}_{(r, d-r-1,1)}} \operatorname{det}\left(G_{A, B, C}^{\mathbf{I}}\right) \tag{6.2}
\end{equation*}
$$

In the following two lemmas, we present some technical results on determinant expansions.

Lemma 6.1. For any $m \geq 1$ and $A_{1}, \ldots, A_{m} \in \mathcal{M}$ we have

$$
\begin{equation*}
\operatorname{det}\left(A_{1}+\cdots+A_{m}\right)=\sum_{\mathbf{p} \in \mathcal{P}_{m}} \sum_{\mathbf{I} \in \mathcal{I}_{\mathbf{p}}} \operatorname{det}\left(G_{A_{1}, \ldots, A_{m}}^{\mathbf{I}}\right) \tag{6.3}
\end{equation*}
$$

Proof. Letting $\mathcal{S}_{d}$ be the set of all permutations of $\{1, \ldots, d\}$ and $\operatorname{sign}(s)$ be the signature of $s \in \mathcal{S}$, we have

$$
\begin{aligned}
\operatorname{det}(A+B) & =\sum_{s \in \mathcal{S}_{d}}(-1)^{\operatorname{sign}(s)} \prod_{i=1}^{d}\left(a^{s(i), i}+b^{s(i), i}\right) \\
& =\sum_{I \subset\{1, \ldots, d\}} \sum_{s \in \mathcal{S}_{d}}(-1)^{\operatorname{sign}(s)} \prod_{i \in I} a^{s(i), i} \prod_{i \notin I} b^{s(i), i} \\
& =\sum_{I \subset\{1, \ldots, d\}} \operatorname{det}\left(G_{A, B}^{\left(I, I^{c}\right)}\right) .
\end{aligned}
$$

This readily implies that if (6.3) holds for some $m$, it also holds for $m+1$. Since (6.3) is obvious for $m=1$, the result follows by induction on $m$.

Lemma 6.2. There is a constant $K$ such that, for all $r=0, \ldots, d$, all $h \in$ $(0,1]$ and all $A, B, C, D \in \mathcal{M}$ with $\operatorname{rank}(A) \leq r$ we have, with $\Lambda=\|A\|+\|B\|+$ $\|C\|+\|D\|$ and with the convention $\gamma_{-1}(A, B)=0$ :

$$
\begin{align*}
& \mid \operatorname{det}\left(A+h B+h^{2} C+h^{2} D\right)-h^{d-r} \gamma_{r}(A, B) \\
& -h^{d-r+1}\left(\gamma_{r-1}(A, B)+\gamma_{r}^{\prime}(A, B, C)\right) \mid  \tag{6.4}\\
& \leq K h^{r-d+1} \Lambda^{d-1}(h \Lambda+\|D\|) \text {, } \\
& \left\lvert\, \frac{1}{h^{2 d-2 r}} \operatorname{det}\left(A+h B+h^{2} C+h^{2} D\right)^{2}-\gamma_{r}(A, B)^{2}\right. \\
& -2 h \gamma_{r}(A, B)\left(\gamma_{r-1}(A, B)+\gamma_{r}^{\prime}(A, B, C)\right)  \tag{6.5}\\
& \leq K h \Lambda^{2 d-1}(h \Lambda+\|D\|) \text {. }
\end{align*}
$$

Proof. Let $\mathbf{p} \in \mathcal{P}_{4}$ and $\mathbf{I} \in \mathcal{I}_{\mathbf{p}}$. Then $\operatorname{det}\left(G_{A, h B, h^{2} C, h^{2} D}^{\mathbf{I}}\right)=h^{p_{2}+2 p_{3}+2 p_{4}} \times$ $\operatorname{det}\left(G_{A, B, C, D}^{\mathbf{I}}\right)$ vanishes when $p_{1}>r$, and has absolute value smaller than $K h^{p_{2}+2 p_{3}+2 p_{4}} \Lambda^{d-p_{4}}\|D\|^{p_{4}}$. Then (6.4) readily follows from (6.3), and by taking squares in (6.4) we deduce (6.5).

With the same notation, and if further $A^{\prime}, B^{\prime}, C^{\prime}, D^{\prime} \in \mathcal{M}$ with $\operatorname{rank}\left(A^{\prime}\right) \leq r$ also and $\Lambda^{\prime}=\left\|A^{\prime}\right\|+\left\|B^{\prime}\right\|+\left\|C^{\prime}\right\|+\left\|D^{\prime}\right\|$, and $h^{\prime} \in(0,1]$, the same argument shows that

$$
\left\lvert\, \frac{1}{\left(h h^{\prime}\right)^{2 d-2 r}} \operatorname{det}\left(A+h B+h^{2} C+h^{2} D\right)^{2}\right.
$$

$$
\begin{align*}
& \times \operatorname{det}\left(A^{\prime}+h^{\prime} B^{\prime}+h^{\prime 2} C^{\prime}+h^{\prime 2} D^{\prime}\right)^{2}-\gamma_{r}(A, B)^{2} \gamma_{r}\left(A^{\prime}, B^{\prime}\right)^{2}  \tag{6.6}\\
& \quad \leq K\left(h+h^{\prime}\right)\left(\Lambda \Lambda^{\prime}\right)^{2 d}
\end{align*}
$$

### 6.2. Proof of Lemma 3.1.

(1) The results about $\Gamma_{r}(\underline{u})$. We write $V_{i}$ and $\bar{V}_{i}$ for the $d$-dimensional variables whose components are, respectively, the $d$ first and the $d$ last components of $\Psi(\underline{u}, 1)_{i}$, for which we can take $\bar{W}=W$ and $\bar{W}^{\prime}=W^{\prime}$, and we set $A=\operatorname{mat}\left(V_{1}, \ldots, V_{d}\right)$ and $B=\operatorname{mat}\left(\bar{V}_{1}, \ldots, \bar{V}_{d}\right)$. If $\Delta_{j} W^{(\prime)}=W_{j}^{(\prime)}-W_{j-1}^{(\prime)}$, we have

$$
\begin{align*}
& V_{i}^{l}=\sum_{m=1}^{q} \alpha^{l m} \Delta_{i} W^{m}  \tag{6.7}\\
& \bar{V}_{i}^{l}=a^{l}+\sum_{m=1}^{d} \beta^{l m} \Delta_{i} W^{\prime m}+\sum_{m, k=1}^{q} \gamma^{l k m} h_{i, k m}(W)
\end{align*}
$$

where each $h_{i, l m}$ is a function of the path of $W$. Note also that $\bar{F}_{r}(\underline{u}, 1)=$ $\gamma_{r}(A, B)^{2}$.

Assuming first that the rank of $\alpha$ is (strictly) smaller than $r$, we observe that the rank of $A$ is also smaller than $r$, implying by (6.1) that $\gamma_{r}(A, B)=0$, hence $\Gamma_{r}(\underline{u})=0$.

Next we assume that the rank of $\alpha$ is $r$, and proceed to prove $\Gamma_{r}(\underline{u})>0$. We first simplify the problem as follows. The matrix $\beta$ is invertible and the rank of $\beta^{-1} \alpha \alpha^{*} \beta^{-1, *}$ is $r$, so we can write $\beta^{-1} \alpha=\Pi \Lambda$, where $\Pi \in \mathcal{M}$ is an orthonormal matrix and $\Lambda \in \mathcal{M}$ is a diagonal matrix whose diagonal entries $\lambda_{j}$ satisfy $\lambda_{j} \neq 0$ if $j \leq r$ and $\lambda_{j}=0$ otherwise. Then, setting $V_{j}^{\prime}=\Pi^{*} \beta^{-1} V_{j}$ and $\bar{V}_{j}^{\prime}=\Pi^{*} \beta^{-1} \bar{V}_{j}$, the sequence $\left(V_{j}^{\prime}, \bar{V}_{j}^{\prime}\right)$ has the form (6.7), upon replacing $W^{\prime}$ by $\Pi^{*} W^{\prime}$ (another Brownian motion) and $\underline{u}=(\alpha, \beta, \gamma, a)$ by $\underline{u}^{\prime}=\left(\alpha^{\prime}, J, \gamma^{\prime}, a^{\prime}\right)$, where $J$ is the identity in $\mathcal{M}$ and $\alpha^{\prime}=\Pi^{*} \beta^{-1} \alpha=$ $\Lambda$ and $\gamma^{\prime i j l}=\sum_{m=1}^{d}\left(\Pi^{*} \beta^{-1}\right)^{j m} \gamma^{m j l}$ and $a^{\prime}=\Pi^{*} \underline{\beta}^{-1} a$. Furthermore, $A^{\prime}=$ $\operatorname{mat}\left(V_{1}^{\prime}, \ldots, V_{d}^{\prime}\right)=\Pi^{*} \beta^{-1} A$ and $B^{\prime}=\operatorname{mat}\left(\bar{V}_{1}^{\prime}, \ldots, \bar{V}_{d}^{\prime}\right)=\Pi^{*} \beta^{-1} B$, implying $\gamma_{r}\left(A^{\prime}, B^{\prime}\right)=\operatorname{det}\left(\Pi^{*} \beta^{-1}\right) \gamma_{r}(A, B)$, which in turn yields $\Gamma_{r}(\underline{u})=\frac{1}{\operatorname{det}(\beta)^{2}} \Gamma_{r}\left(\underline{u}^{\prime}\right)$, because $\operatorname{det}(\Pi)=1$.

In other words, it is enough to prove the result when $\beta=J$ and $\alpha=\Lambda$ is diagonal as above, and below we assume this. The two matrices $A$ and $B$ can thus be realized as

$$
A=\left(\begin{array}{ccc}
\lambda_{1} \Phi_{1}^{1} & \cdots & \lambda_{1} \Phi_{d}^{1}  \tag{6.8}\\
\cdot & \cdots & \cdot \\
\cdot & \cdots & \cdot \\
\lambda_{d} \Phi_{1}^{d} & \cdots & \lambda_{d} \Phi_{d}^{d}
\end{array}\right), \quad B=\left(\begin{array}{ccc}
\Upsilon_{1}^{1}+\Theta_{1}^{1} & \ldots & \Upsilon_{d}^{1}+\Theta_{d}^{1} \\
\cdot & \cdots & \cdot \\
\cdot & \cdots & \cdot \\
\Upsilon_{1}^{d}+\Theta_{1}^{d} & \cdots & \Upsilon_{d}^{d}+\Theta_{d}^{d}
\end{array}\right)
$$

where all $\Phi_{j}^{i}$ and $\Upsilon_{j}^{i}$ are i.i.d. $\mathcal{N}(0,1)$ and the variables $\Theta_{j}^{i}$ are independent of the $\Upsilon_{m}^{l}$ 's (note that we have incorporated the constant $a^{i}$ in each variable $\Theta_{j}^{i}$ ).

Let $\mathcal{J}_{r}$ be the class of all subsets of $\{1, \ldots, d\}$ with $r$ points. Since $\lambda_{j} \neq 0$ if $j \leq r$ and $\lambda_{j}=0$ otherwise, we see that, if $\mathbf{I}=\left(I, I^{c}\right)$ with $I=\left\{j_{1}<\cdots<j_{r}\right\} \in$
$\mathcal{J}_{r}$ and $I^{c}=\left\{j_{1}^{\prime}<\cdots<j_{d-r}^{\prime}\right\}$, we have $\operatorname{det}\left(G_{A, B}^{\mathbf{I}}\right)=\varepsilon_{I} \operatorname{det}\left(A_{I}\right) \operatorname{det}\left(B_{I}\right)$, where $A_{I}$ and $B_{I}$ are the $r \times r$ and $(d-r) \times(d-r)$ matrices with entries $A_{I}^{l, m}=$ $A^{j_{l}, m}$ and $B_{I}^{l, m}=B^{j_{l}^{\prime}, r+l}$, and $\varepsilon_{I}$ takes values in $\{-1,1\}$. Thus,

$$
\gamma_{r}(A, B)=\sum_{I \in \mathcal{J}_{r}} \varepsilon_{I} \operatorname{det}\left(A_{I}\right) \operatorname{det}\left(B_{I}\right)
$$

In this sum, we single out the $I$ 's which contain $d$, and those which do not, and for the former ones the product $\operatorname{det}\left(A_{I}\right) \operatorname{det}\left(B_{I}\right)$ does not depend of the vector $\Upsilon_{d}$. For those which do not contain $d$, we develop $\operatorname{det}\left(B_{I}\right)$ along the last column, which involves the determinants of the matrices $B_{I, i}$ which are the restrictions of $B$ to the last $d-r$ lines except $i$, and to the column indexed by the complement $I^{c}$ of $I$, except $d$. We thus get

$$
\begin{align*}
\gamma_{r}(A, B) & =Z+\sum_{i=1}^{d-r}(-1)^{i}\left(\Upsilon_{d}^{r+i}+\Theta_{d}^{r+i}\right) Z_{i}^{\prime},  \tag{6.9}\\
Z_{i}^{\prime} & =\sum_{I \in \mathcal{J}_{r}, d \notin I} \varepsilon_{I} \operatorname{det}\left(A_{I}\right) \operatorname{det}\left(B_{I, i}\right),
\end{align*}
$$

where $Z$ and all $Z_{i}^{\prime}$ and $\Theta_{d}^{r+i}$ are independent of the vector $\Upsilon_{d}$. Since this random vector $\Upsilon_{d}$ has a density, it follows that the variable $\gamma_{r}(A, B)$ also has a density, provided $Z_{i}^{\prime} \neq 0$ a.s. for at least one value of $i$.

At this stage, we observe that $Z_{i}^{\prime}$ has exactly the same structure as $\gamma_{r}(A, B)$, except that the dimension of each $B_{I, i}$ is $(d-r-1) \times(d-r-1)$ instead of $(d-r) \times(d-r)$, and that the last column of the original problem has totally disappeared. We can repeat the argument, to obtain that $Z_{i}^{\prime}$ has a density and is thus a.s. nonvanishing, as soon as some similar quantity (where the last two columns of the original problem no longer show up) is a.s. nonvanishing. Then, after an obvious induction, we deduce that $\gamma_{r}(A, B)$ has a density as soon as $\operatorname{det}\left(A_{I}\right) \neq 0$ a.s. for $I=\{1, \ldots, r\}$.

However, since the entries of this last $A_{I}$ are $\lambda_{i} \Phi_{j}^{i}$ for $i, j=1, \ldots, r$, and all those $\lambda_{i}$ are nonzero, it is well known [and also a simple consequence of the previous proof, in which we develop $\operatorname{det}\left(A_{I}\right)$ according to its last column and perform the same induction procedure] that $\operatorname{det}\left(A_{I}\right)$ has a density. This indeed shows us that $\gamma_{r}(A, B)$ has a density, hence $\mathbb{E}\left(\gamma_{r}(A, B)^{2}\right)>0$ and the proof of the first part of (3.5) is complete.
(2) The results about $\Gamma_{r}^{\prime}(\underline{u})$ and $\Gamma_{r}^{\prime \prime}(\underline{u})$. When the rank of $\alpha$ is smaller than $r$, we have seen that, with the previous notation, $\gamma_{r}(A, B)=0$, hence also $\Gamma_{r}^{\prime}(\underline{u})=$ $\Gamma_{r}^{\prime \prime}(\underline{u})=0$.

Next, we turn to the case when the rank of $\alpha$ is $r$. Exactly as in the previous proof, it suffices to show the result when $\underline{u}=(\Lambda, J, \gamma, a)$. Recalling that $\bar{F}_{r}(\underline{u}, 1)$ and $\bar{F}_{r}(\underline{u}, 2)$ have the same law, we have $\overline{\mathbb{E}}\left(\left(\bar{F}_{r}(\underline{u}, 1)-\bar{F}_{r}(\underline{u}, 2)\right)^{2}\right)=$
$2\left(\Gamma_{r}^{\prime}(\underline{u})-\Gamma_{r}^{\prime \prime}(\underline{u})\right)$ and thus the second part of (3.5) holds unless $\bar{F}_{r}(\underline{u}, 1)=$ $\bar{F}_{r}(u, 2)$ a.s.

With the previous notation, we have $\bar{F}_{r}(\underline{u}, 1)=\gamma_{r}(A, B)^{2}$, and also $\bar{F}_{r}(\underline{u}, 2)=$ $\gamma_{r}(\bar{A}, \bar{B})^{2}$, where $\bar{A}$ and $\bar{B}$ are given again by (6.8) with the same $\lambda_{j}$ 's, and random vectors ( $\bar{\Phi}_{j}, \bar{\Upsilon}_{j}, \bar{\Theta}_{j}$ ) having globally the same distribution as ( $\Phi_{j}, \Upsilon_{j}, \Theta_{j}$ ) (which may of course be defined for $j>d$ ): these two families of vectors are not independent, and we have in fact

$$
\begin{equation*}
\bar{\Phi}_{j}=\frac{1}{\sqrt{2}}\left(\Phi_{2 j-1}+\Phi_{2 j}\right), \quad \bar{\Upsilon}_{j}=\frac{1}{\sqrt{2}}\left(\Upsilon_{2 j-1}+\Upsilon_{2 j}\right) \tag{6.10}
\end{equation*}
$$

plus a more complicated relation relating $\bar{\Theta}_{j}$ with the $\Theta_{j^{\prime}}$ for $j^{\prime} \leq 2 j$ and the vector $a$. What we need to prove is then $\mathbb{P}\left(\left|\gamma_{r}(A, B)\right| \neq\left|\gamma_{r}(\bar{A}, \bar{B})\right|\right)>0$.

We have (6.9), and also, by the same argument,

$$
\begin{aligned}
\gamma_{r}(\bar{A}, \bar{B}) & =\bar{Z}+\sum_{i=1}^{d-r}(-1)^{i}\left(\frac{\Upsilon_{2 d+1}^{r+i}+\Upsilon_{2 d}^{r+i}}{\sqrt{2}}+\bar{\Theta}_{d}^{r+i}\right) \bar{Z}_{i}^{\prime} \\
\bar{Z}_{i}^{\prime} & =\sum_{I \in \mathcal{J}_{r}, d \notin I} \varepsilon_{I} \operatorname{det}\left(\bar{A}_{I}\right) \operatorname{det}\left(\bar{B}_{I, i}\right),
\end{aligned}
$$

where we have also used the second part of (6.10). Here, the vector $\Upsilon_{2 d}$ has a density and is independent of all other terms showing in the above expression, and also independent of $\gamma_{r}(A, B)$. Therefore, $\left|\gamma_{r}(A, B)\right| \neq\left|\gamma_{r}(\bar{A}, \bar{B})\right|$ almost surely on the set $\left\{\bar{Z}_{i}^{\prime} \neq 0\right\}$. Now, $\bar{Z}_{i}^{\prime}$ is the same as $Z_{i}^{\prime}$, upon replacing $(A, B)$ by $(\bar{A}, \bar{B})$ everywhere, hence the previous proof shows that indeed $\bar{Z}_{i}^{\prime} \neq 0$ a.s. This shows that in fact $\mathbb{P}\left(\left|\gamma_{r}(A, B)\right| \neq\left|\gamma_{r}(\bar{A}, \bar{B})\right|\right)=1$, thus ending the proof of the second part of (3.5).
6.3. Some stochastic calculus preliminaries. We assume Assumption (H) and, by localization (see, e.g., Section 3 in [1]), we may also assume that all processes $X_{t}, \sigma_{t}, b_{t}, a_{t}, v_{t}, a_{t}^{\prime}, v_{t}^{\prime}, a_{t}^{\prime \prime}, v_{t}^{\prime \prime}$ are uniformly bounded in $(\omega, t)$. The constants are always written as $K$, or $K_{p}$ if we want to stress the dependency on an additional parameter $p$, and never depend on $t, i, n, j$. For any process $Y$, we use the following simplifying notation:

$$
\begin{equation*}
\mathcal{F}_{i}^{n}=\mathcal{F}_{2 i d \Delta_{n}}, \quad Y_{i}^{n}=Y_{2 i d \Delta_{n}} \tag{6.11}
\end{equation*}
$$

For all $p, t, s>0$, we have by the Burkholder-Gundy inequality

$$
\begin{equation*}
\mathbb{E}\left(\sup _{u \in[0, s]}\left|Y_{t+u}-Y_{t}\right|^{p} \mid \mathcal{F}_{t}\right) \leq K_{p} s^{p / 2} \quad \text { if } Y=X, \sigma, b, v \tag{6.12}
\end{equation*}
$$

We set

$$
\begin{align*}
\eta_{t, s} & =\sup _{u \in[0, s], Y=a, v^{\prime}, v^{\prime \prime}}\left\|Y_{t+u}-Y_{t}\right\|^{2} \\
\eta_{i}^{n} & =\sqrt{\mathbb{E}\left(\eta_{2 i d \Delta_{n}, 2 d \Delta_{n}} \mid \mathcal{F}_{i}^{n}\right)} . \tag{6.13}
\end{align*}
$$

LEMMA 6.3. For all $t>0$, we have $\Delta_{n} \mathbb{E}\left(\sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \eta_{i}^{n}\right) \rightarrow 0$.
Proof. It suffices to prove the result separately when $Y=a$ or $Y=v^{\prime}$ or $Y=v^{\prime \prime}$. Set $\gamma_{t}^{n}=\sup _{s \in\left(0,4 d \Delta_{n}\right]}\left\|Y_{t+s}-Y_{t}\right\|^{2}$, so $\mathbb{E}\left(\left(\eta_{i}^{n}\right)^{2}\right)$ is smaller than $\mathbb{E}\left(\gamma_{0}^{n}\right)$ when $i=0$ and than $\frac{1}{2 d \Delta_{n}} \int_{2(i-1) d \Delta_{n}}^{2 i d \Delta_{n}} \mathbb{E}\left(\gamma_{s}^{n}\right) d s$ when $i \geq 1$. Hence by the CauchySchwarz inequality,

$$
\begin{aligned}
\Delta_{n} \mathbb{E}\left(\sum_{i=1}^{\left[t / 2 d \Delta_{n}\right]-1} \eta_{i}^{n}\right) & \leq \frac{\sqrt{t}}{\sqrt{2 d}}\left(\mathbb{E}\left(\Delta_{n} \sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1}\left(\eta_{i}^{n}\right)^{2}\right)\right)^{1 / 2} \\
& \leq K_{t}\left(\mathbb{E}\left(\gamma_{0}^{n}+\int_{0}^{t} \gamma_{s}^{n} d s\right)\right)^{1 / 2}
\end{aligned}
$$

We have $\gamma_{s}^{n} \leq K$, whereas the càdlàg property of $Y$ yields that $\gamma_{s}^{n}(\omega) \rightarrow 0$ for all $\omega$, and all $s$ except for countably many strictly positive values (depending on $\omega$ ). Then, the claim follows by the dominated convergence theorem.

The proof of Theorem 3.2 is based on a decomposition of the increments $Z_{(2 i d+\kappa j) \Delta_{n}}^{n, \kappa}-Z_{(2 i d+\kappa(j-1)) \Delta_{n}}^{n, \kappa}$. In order to understand better this decomposition, we first deduce from (2.8) that, for any $z \leq t \leq s$, and with vector notation,

$$
\begin{aligned}
\int_{t}^{s} b_{u} d u & =\alpha_{1}+\alpha_{2}+\alpha_{3}+\alpha_{4} \\
\int_{t}^{s} \sigma_{u} d W_{u} & =\alpha_{5}+\alpha_{6}+\alpha_{7}+\alpha_{8}+\alpha_{9}+\alpha_{10}+\alpha_{11}
\end{aligned}
$$

where

$$
\begin{aligned}
& \alpha_{1}=b_{z}(s-t), \quad \alpha_{2}=\int_{t}^{s}\left(\int_{z}^{u} a_{w}^{\prime} d w\right) d u \\
& \alpha_{3}=v_{z}^{\prime} \int_{t}^{s}\left(W_{u}-W_{z}\right) d u, \quad \alpha_{4}=\int_{t}^{s}\left(\int_{z}^{u}\left(v_{w}^{\prime}-v_{z}^{\prime}\right) d W_{w}\right) d u \\
& \alpha_{5}=\sigma_{z}\left(W_{s}-W_{t}\right), \quad \alpha_{6}=a_{z} \int_{t}^{s}(u-z) d W_{u} \\
& \alpha_{7}=\int_{t}^{s}\left(\int_{z}^{u}\left(a_{w}-a_{z}\right) d w\right) d W_{u} \\
& \alpha_{8}=v_{z} \int_{t}^{s}\left(W_{u}-W_{z}\right) d W_{u}, \quad \alpha_{9}=\int_{t}^{s}\left(\int_{z}^{u}\left(\int_{z}^{w} a_{r}^{\prime \prime} d r\right) d W_{w}\right) d W_{u} \\
& \alpha_{10}=v_{z}^{\prime \prime} \int_{t}^{s}\left(\int_{z}^{u}\left(W_{w}-W_{z}\right) d W_{w}\right) d W_{u} \\
& \alpha_{11}=\int_{t}^{s}\left(\int_{z}^{u}\left(\int_{z}^{w}\left(v_{r}^{\prime \prime}-v_{z}^{\prime \prime}\right) d W_{r}\right) d W_{w}\right) d W_{u} .
\end{aligned}
$$

A repeated use of the Burkholder-Gundy and Hölder inequalities shows that, in view of our assumptions on the various coefficients, and for any $p \geq 1$ :

$$
\mathbb{E}\left(\left\|\alpha_{j}\right\|^{p} \mid \mathcal{F}_{z}\right) \leq \begin{cases}K_{p}(s-z)^{p / 2}, & \text { if } j=5, \\ K_{p}(s-z)^{p}, & \text { if } j=1,8 \\ K_{p}(s-z)^{3 p / 2}, & \text { if } j=3,6,10 \\ K_{p}(s-z)^{2 p}, & \text { if } j=2,9 \\ K_{p}(s-z)^{3 p / 2} \mathbb{E}\left(\eta_{z, s-z}^{p} \mid \mathcal{F}_{z}\right), & \text { if } j=4,7,11\end{cases}
$$

We can then apply the previous decomposition with $z=2 i d \Delta_{n}, t=(2 i d+$ $\kappa(j-1)) \Delta_{n}$ and $s=(2 i d+\kappa j) \Delta_{n}$, and add the increment of the process $X^{\prime}$, to obtain

$$
\begin{align*}
& \frac{Z_{(2 i d+j \kappa) \Delta_{n}}^{n, \kappa}-Z_{(2 i d+(j-1) \kappa) \Delta_{n}}^{n, \kappa}}{\sqrt{\kappa \Delta_{n}}}  \tag{6.14}\\
& \quad=\alpha_{i, j}^{n, \kappa}+\sqrt{\kappa \Delta_{n}} \beta_{i, j}^{n, \kappa}+\kappa \Delta_{n} \gamma_{i, j}^{n, \kappa}+\Delta_{n} \delta_{i, j}^{n, \kappa}
\end{align*}
$$

for $\kappa=1,2$, and where (explicitly writing the components)

$$
\begin{aligned}
\alpha_{i, j}^{n, \kappa, l}= & \frac{1}{\sqrt{\kappa \Delta_{n}}} \sum_{m=1}^{q} \sigma_{i}^{n, l m}\left(W_{(2 i d+\kappa j) \Delta_{n}}^{m}-W_{(2 i d+\kappa(j-1)) \Delta_{n}}^{m}\right), \\
\beta_{i, j}^{n, \kappa, l}= & b_{i}^{n, l}+\frac{1}{\kappa \Delta_{n}} \sum_{m, k=1}^{q} v_{i}^{n, l m k} \int_{(2 i d+\kappa(j-1)) \Delta_{n}}^{(2 i d+\kappa j) \Delta_{n}}\left(W_{s}^{k}-W_{2 i d \Delta_{n}}^{k}\right) d W_{s}^{m} \\
& +\frac{1}{\sqrt{\kappa \Delta_{n}}} \sum_{m=1}^{d} \tilde{\sigma}^{l m}\left(W_{(2 i d+\kappa j) \Delta_{n}}^{\prime m}-W_{(2 i d+\kappa(j-1)) \Delta_{n}}^{\prime m}\right), \\
\gamma_{i, j}^{n, \kappa, l=}= & \frac{1}{\left(\kappa \Delta_{n}\right)^{3 / 2}} \\
& \times\left(\sum_{m=1}^{q} a_{i}^{n, l m} \int_{(2 i d+\kappa(j-1)) \Delta_{n}}^{(2 i d+\kappa j) \Delta_{n}}\left(s-2 i d \Delta_{n}\right) d W_{s}^{m}\right. \\
& +\sum_{m=1}^{q} v_{i}^{\prime n, l m} \int_{(2 i d+\kappa(j-1)) \Delta_{n}}^{(2 i d+\kappa j) \Delta_{n}}\left(W_{s}^{m}-W_{2 i d \Delta_{n}}^{m}\right) d s \\
& +\sum_{m, k=1}^{q} v_{i}^{\prime \prime n, m l k} \\
& \left.\times \int_{(2 i d+\kappa(j-1)) \Delta_{n}}^{(2 i d+\kappa j) \Delta_{n}}\left(\int_{(2 i d+\kappa(j-1)) \Delta_{n}}^{s}\left(W_{u}^{k}-W_{2 i d \Delta_{n}}^{k}\right) d W_{u}^{l}\right) d W_{s}^{m}\right)
\end{aligned}
$$

and $\delta_{i, j}^{n, \kappa}$ is a remainder term, and for $p \geq 1$ we have the estimates when $j \leq 2 d$ if $\kappa=1$ and $j \leq d$ when $\kappa=2$ (recalling $\eta_{t, s} \leq K$ ):

$$
\begin{align*}
& \mathbb{E}\left(\left\|\alpha_{i, j}^{n, \kappa}\right\|^{p}+\left\|\beta_{i, j}^{n, \kappa}\right\|^{p}+\left\|\gamma_{i, j}^{n, \kappa}\right\|^{p} \mid \mathcal{F}_{i}^{n}\right) \leq K_{p}  \tag{6.15}\\
& \mathbb{E}\left(\left\|\delta_{i, j}^{n, \kappa}\right\|^{p} \mid \mathcal{F}_{i}^{n}\right) \leq K_{p}\left(\Delta_{n}^{p / 2}+\left(\eta_{i}^{n}\right)^{2 \wedge p}\right) \leq K_{p}
\end{align*}
$$

We end these preliminaries with a lemma which compares $S^{n, \kappa}$ for $\kappa=1,2$ with the following processes:

$$
S(r)_{t}^{n, \kappa}=2 d \Delta_{n} \sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \gamma_{r}\left(A_{i}^{n, \kappa}, B_{i}^{n, \kappa}\right)^{2}
$$

$$
\begin{equation*}
\text { where } A_{i}^{n, \kappa}=\operatorname{mat}\left(\alpha_{i, 1}^{n, \kappa}, \ldots, \alpha_{i, d}^{n, \kappa}\right), B_{i}^{n, \kappa}=\operatorname{mat}\left(\beta_{i, 1}^{n, \kappa}, \ldots, \beta_{i, d}^{n, \kappa}\right) \tag{6.16}
\end{equation*}
$$

It also compares $V^{n, \kappa, \kappa^{\prime}}$ of (3.15) with

$$
\begin{equation*}
V(r)_{t}^{n, \kappa, \kappa^{\prime}}=4 d^{2} \Delta_{n} \sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \gamma_{r}\left(A_{i}^{n, \kappa}, B_{i}^{n, \kappa}\right)^{2} \gamma_{r}\left(A_{i}^{n, \kappa^{\prime}}, B_{i}^{n, \kappa^{\prime}}\right)^{2} \tag{6.17}
\end{equation*}
$$

LEMMA 6.4. If $r_{t}(\omega) \leq r$ identically for some $r \in\{0, \ldots, d\}$, we have for $\kappa$, $\kappa^{\prime}=1,2$ :

$$
\begin{equation*}
\frac{1}{\sqrt{\Delta_{n}}}\left(\frac{1}{\left(\kappa \Delta_{n}\right)^{d-r}} S^{n, \kappa}-S(r)^{n, \kappa}\right) \stackrel{\text { u.c.p. }}{\Longrightarrow} 0 \tag{6.18}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{1}{\left(\kappa \kappa^{\prime} \Delta_{n}^{2}\right)^{d-r}} V^{n, \kappa, \kappa^{\prime}}-V(r)^{n, \kappa, \kappa^{\prime}} \stackrel{\text { u.c.p. }}{\Longrightarrow} 0 \tag{6.19}
\end{equation*}
$$

Proof. We denote by $\xi_{i}^{n, \kappa}$ the $i$ th summand in the definition (2.13) of $S_{t}^{n, \kappa}$. Besides the matrices in (6.16), we also define

$$
C_{i}^{n, \kappa}=\operatorname{mat}\left(\gamma_{i, 1}^{n, \kappa}, \ldots, \gamma_{i, d}^{n, \kappa}\right), \quad D_{i}^{n, \kappa}=\operatorname{mat}\left(\delta_{i, 1}^{n, \kappa}, \ldots, \delta_{i, d}^{n, \kappa}\right)
$$

We start with (6.18). Applying (6.5) with $h=\sqrt{\kappa \Delta_{n}}$, the fact that each $A_{i}^{n, \kappa}$ has at most rank $r$ (because $r_{t} \leq r$ ), and the estimates (6.15) plus the Cauchy-Schwarz inequality, we obtain

$$
\frac{1}{\left(\kappa \Delta_{n}\right)^{d-r}} \xi_{i}^{n, \kappa}=\gamma_{r}\left(A_{i}^{n, \kappa}, B_{i}^{n, \kappa}\right)^{2}+2 \sqrt{\kappa \Delta_{n}} \zeta_{i}^{n, \kappa}+\widetilde{\zeta}_{i}^{n, \kappa},
$$

where

$$
\begin{aligned}
\zeta_{i}^{n, \kappa} & =\gamma_{r}\left(A_{i}^{n, \kappa}, B_{i}^{n, \kappa}\right)\left(\gamma_{r-1}\left(A_{i}^{n, \kappa}, B_{i}^{n, \kappa}\right)+\gamma_{r}^{\prime}\left(A_{i}^{n, \kappa}, B_{i}^{n, \kappa}, C_{i}^{n, \kappa}\right)\right) \\
\mathbb{E}\left(\left|\tilde{\zeta}_{i}^{n, \kappa}\right|\right) & \leq K \Delta_{n}+K \sqrt{\Delta_{n}} \mathbb{E}\left(\eta_{i}^{n}\right)
\end{aligned}
$$

In view of Lemma 6.3, $\sqrt{\Delta_{n}} \sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \widetilde{\zeta}_{i}^{n, \kappa} \xrightarrow{\text { u.c.p. }} 0$. Since $S_{t}^{n, \kappa}=2 d \Delta_{n} \times$ $\sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \xi_{i}^{n, \kappa}$, it remains to prove that $\Delta_{n} \sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \zeta_{i}^{n, \kappa} \xrightarrow{\text { u.c.p. }} 0$.

For this purpose we use the decomposition $\zeta_{i}^{n, \kappa}=\zeta_{i}^{\prime n, \kappa}+\zeta_{i}^{\prime \prime n, \kappa}$, where $\zeta_{i}^{\prime n, \kappa}=\mathbb{E}\left(\zeta_{i}^{n, \kappa} \mid \mathcal{F}_{i}^{n}\right)$. By Doob's inequality, (6.15) and the fact that $\zeta_{i}^{\prime \prime n}$ is $\mathcal{F}_{i+1}^{n}$-measurable, we have

$$
\mathbb{E}\left(\sup _{s \leq t}\left(\sum_{i=0}^{\left[s / 2 d \Delta_{n}\right]-1} \zeta_{i}^{\prime \prime n, \kappa}\right)^{2}\right) \leq 2^{d+1} \mathbb{E}\left(\sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1}\left|\zeta_{i}^{n}\right|^{2}\right) \leq \frac{K t}{\Delta_{n}}
$$

Thus $\Delta_{n} \sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \zeta_{i}^{\prime \prime n, \kappa} \stackrel{\text { u.c.p. }}{\Longrightarrow} 0$, and the result will hold if we can prove that $\zeta_{i}^{\prime n, \kappa}=0$. We even prove the stronger statement that $\mathbb{E}\left(\zeta_{i}^{n, \kappa} \mid \mathcal{G}^{W^{\prime}} \vee \mathcal{F}_{i}^{n}\right)=0$, where $\mathcal{G}^{W^{\prime}}$ is the $\sigma$-field generated by the whole process $W^{\prime}$, and this is implied by

$$
\begin{align*}
& \mathbf{I} \in \mathcal{I}_{(r, d-r)}, \mathbf{I}^{\prime} \in \mathcal{I}_{(r-1, d-r+1)}, \mathbf{I}^{\prime \prime} \in \mathcal{I}_{(r, d-r-1,1)} \\
& \Longrightarrow \mathbb{E}\left(\operatorname{det}\left(G_{A_{i}^{n, \kappa}, B_{i}^{n, \kappa}}^{\mathbf{I}}\right) \operatorname{det}\left(G_{A_{i}^{n, \kappa}, B_{i}^{n, \kappa}}^{\mathbf{I}^{\prime}}\right) \mid \mathcal{G}^{W^{\prime}} \vee \mathcal{F}_{i}^{n}\right)=0,  \tag{6.20}\\
& \mathbb{E}\left(\operatorname{det}\left(G_{A_{i}^{n, \kappa}, B_{i}^{n, \kappa}}^{\mathbf{I}}\right) \operatorname{det}\left(G_{A_{i}^{n, \kappa}, B_{i}^{n, \kappa}, C_{i}^{n, \kappa}}^{\mathbf{I}^{\prime \prime}}\right) \mid \mathcal{G}^{W^{\prime}} \vee \mathcal{F}_{i}^{n}\right)=0 .
\end{align*}
$$

The variables $\alpha_{i, j}^{n, \kappa, l}, \beta_{i, j}^{n, \kappa, l}$ and $\gamma_{i, j}^{n, \kappa, l}$ have the form $\Phi\left(\omega,\left(W(\omega)_{2 i d \Delta_{n}+t}-\right.\right.$ $\left.W(\omega)_{2 i d \Delta_{n}}\right)_{t \geq 0}$, with $\Phi$ a $\left(\mathcal{G}^{W^{\prime}} \vee \mathcal{F}_{i}^{n}\right) \otimes \mathcal{C}^{q}$-measurable function on $\Omega \times$ $C\left(\mathbb{R}_{+}, \mathbb{R}^{q}\right)$, where $C\left(\mathbb{R}_{+}, \mathbb{R}^{q}\right)$ is the set of all continuous $\mathbb{R}^{q}$-valued functions on $\mathbb{R}_{+}$and $\mathcal{C}^{q}$ is its Borel $\sigma$-field for the local uniform topology. When $\Phi=\alpha_{i, j}^{n, \kappa, l}$ or $\Phi=\gamma_{i, j}^{n, \kappa, l}$, the map $x \mapsto \Phi(\omega, x)$ is odd, in the sense that $\Phi(\omega,-x)=-\Phi(\omega, x)$, and it is even when $\Phi=\beta_{i, j}^{n, \kappa, l}$.

In (6.20), the three variables $\operatorname{det}\left(G_{A_{i}^{n, \kappa}, B_{i}^{n, \kappa}}^{\mathbf{I}}\right), \operatorname{det}\left(G_{A_{i}^{n, \kappa}, B_{i}^{n, \kappa}}^{\mathbf{I}^{\prime}}\right), \operatorname{det}\left(G_{\left.A_{i}^{n, \kappa}, B_{i}^{n, \kappa}, C_{i}^{n, \kappa}\right)}^{\mathbf{I}^{\prime \prime}}\right.$ are associated with three functions $\Psi, \Psi^{\prime}, \Psi^{\prime \prime}$ of the same type. What precedes yields that $\Psi$ is even (resp., odd) if $r$ is even (resp., odd), and both $\Psi^{\prime}$ and $\Psi^{\prime \prime}$ are even (resp., odd) if $r$ is odd (resp., even). Consequently, the products $\Psi \Psi^{\prime}$ and $\Psi \Psi^{\prime \prime}$ are odd in all cases. Since the $\mathcal{G}^{W^{\prime}} \vee \mathcal{F}_{i}^{n}$-conditional law of $\left(W_{2 i d \Delta_{n}+t}-W_{2 i d \Delta_{n}}\right)_{t \geq 0}$ is invariant by the map $x \mapsto-x$ on $C\left(\mathbb{R}_{+}, \mathbb{R}^{d}\right)$, we deduce (6.20), hence (6.18) holds.

Finally, we turn to (6.19). Let $\theta_{i}^{n}$ be the $i$ th summand in the right-hand side of (3.15), for $\kappa, \kappa^{\prime}$ fixed. We can apply (6.6) with $h=\sqrt{\kappa \Delta_{n}}$ and $h^{\prime}=\sqrt{\kappa^{\prime} \Delta_{n}}$, and (6.14) and (6.15) again, to get

$$
\mathbb{E}\left(\left|\frac{1}{\left(\kappa \kappa^{\prime} \Delta_{n}^{2}\right)^{d-r}} \theta_{i}^{n}-\gamma_{r}\left(A_{i}^{n, \kappa}, B_{i}^{n, \kappa}\right)^{2} \gamma_{r}\left(A_{i}^{n, \kappa^{\prime}}, B_{i}^{n, \kappa^{\prime}}\right)^{2}\right|\right) \leq K \sqrt{\Delta_{n}} .
$$

Equation (6.19) follows, and the proof is complete.
6.4. Proof of Theorem 3.2, Corollary 3.3 and Proposition 3.4.
(1) Observe that, with the notation (3.2) and (3.3), and upon taking

$$
\begin{align*}
& \underline{u}_{i}^{n}=\left(\sigma_{i}^{n}, \tilde{\sigma}, v_{i}^{n}, b_{i}^{n}\right), \\
& \bar{W}_{t}=\frac{W_{(2 i d+t) \Delta_{n}}-W_{2 i d \Delta_{n}}}{\sqrt{\Delta_{n}}},  \tag{6.21}\\
& \bar{W}_{t}^{\prime}=\frac{W_{(2 i d+t) \Delta_{n}}^{\prime}-W_{2 i d \Delta_{n}}^{\prime}}{\sqrt{\Delta_{n}}},
\end{align*}
$$

we have $\gamma_{r}\left(A_{i}^{n, \kappa}, B_{i}^{n, \kappa}\right)^{2}=\bar{F}\left(\underline{u}_{i}^{n}, \kappa\right)$. We consider the two-dimensional variables $\xi_{i}^{n}$ with components

$$
\begin{equation*}
\xi_{i}^{n, \kappa}=2 d \sqrt{\Delta_{n}}\left(\gamma_{r}\left(A_{i}^{n, \kappa}, B_{i}^{n, \kappa}\right)^{2}-\Gamma_{r}\left(u_{i}^{n}\right)\right), \quad \kappa=1,2 . \tag{6.22}
\end{equation*}
$$

Since $\underline{u}_{i}^{n}$ is $\mathcal{F}_{i}^{n}$-measurable, whereas the processes $\bar{W}$ and $\bar{W}^{\prime}$ above are independent of $\mathcal{F}_{i}^{n}$, we deduce from (3.4), and from (6.15) for the estimate below, that

$$
\begin{align*}
\mathbb{E}\left(\xi_{i}^{n, \kappa} \mid \mathcal{F}_{i}^{n}\right) & =0, \quad \mathbb{E}\left(\left\|\xi_{i}^{n}\right\|^{4} \mid \mathcal{F}_{i}^{n}\right) \leq K \Delta_{n}^{2},  \tag{6.23}\\
\mathbb{E}\left(\xi_{i}^{n, \kappa} \xi_{i}^{n, \kappa^{\prime}} \mid \mathcal{F}_{i}^{n}\right) & = \begin{cases}4 d^{2} \Delta_{n} \Gamma_{r}^{\prime}\left(\underline{u}_{i}^{n}\right), & \text { if } \kappa=\kappa^{\prime}, \\
4 d^{2} \Delta_{n} \Gamma_{r}^{\prime \prime}\left(\underline{u}_{i}^{n}\right), & \text { if } \kappa \neq \kappa^{\prime} .\end{cases}
\end{align*}
$$

(2) By (6.18), for Theorem 3.2 it is enough to prove the stable convergence $U^{\prime}(r)^{n} \stackrel{\mathcal{L}-s}{\Longrightarrow} \mathcal{U}(r)$, where $U^{\prime}(r)^{n}$ is the two-dimensional process with components $U^{\prime}(r)^{n, \kappa}=\frac{1}{\sqrt{\Delta_{n}}}\left(S(r)^{n, \kappa}-S(r)\right)$ and the quantity $S(r)^{n, \kappa}$ is defined in (6.16). We have $U^{\prime}(r)^{n}=Y^{n}+Y^{\prime n}$, where

$$
\begin{aligned}
Y_{t}^{n} & =\sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \xi_{i}^{n}, \\
Y_{t}^{\prime n} & =\frac{1}{\sqrt{\Delta_{n}}}\left(2 d \Delta_{n} \sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \Gamma_{r}\left(\sigma_{i}^{n}, \tilde{\sigma}, v_{i}^{n}, b_{i}^{n}\right)-\int_{0}^{t} \Gamma_{r}\left(\sigma_{s}, \tilde{\sigma}, v_{s}, b_{s}\right) d s\right)
\end{aligned}
$$

and $\xi_{i}^{n}$ is given in (6.22). Since the three processes $\sigma, v, b$ are Itô semimartingales, whereas $\Gamma_{r}$ is a $C^{\infty}$ function, it is well known that $Y^{\prime n} \xrightarrow{\text { u.c.p. }} 0$ (see, e.g., Section 8 in [1]). We are thus left to prove that

$$
\begin{equation*}
Y^{n} \xrightarrow{\mathcal{L}-s} \mathcal{U}(r) \tag{6.24}
\end{equation*}
$$

By virtue of the first two parts of (6.23), a standard CLT for triangular arrays of martingale (see [8], Theorem IX.7.28) increments shows that, for (6.24) to hold, it
suffices to show the next two properties:

$$
\begin{align*}
\sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \mathbb{E}\left(\xi_{i}^{n, \kappa} \xi_{i}^{n, \kappa^{\prime}} \mid \mathcal{F}_{i}^{n}\right) \xrightarrow{\mathbb{P}} V(r)_{t}^{\kappa \kappa^{\prime}},  \tag{6.25}\\
\sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \mathbb{E}\left(\xi_{i}^{n}\left(N_{2(i+1) d \Delta_{n}}-N_{2 i d \Delta_{n}}\right) \mid \mathcal{F}_{i}^{n}\right) \xrightarrow{\mathbb{P}} 0 \tag{6.26}
\end{align*}
$$

for all $t>0$ and for any bounded martingale $N$ orthogonal to ( $W, W^{\prime}$ ) and also for $N=W^{m}$ or $N=W^{\prime m}$ for any $m$.

The last part of (6.23) and the càdlàg property of $\sigma, v, b$, plus the fact that $\Gamma_{r}^{\prime}$ and $\Gamma_{r}^{\prime \prime}$ are polynomials, immediately gives us (6.25) by Riemann integration.

The proof of (6.26) is also standard: By construction, $\xi_{i}^{n}$ is a two-dimensional variable of the form $\Phi\left(\omega,\left(W_{2 i d \Delta_{n}+t}(\omega)-W_{2 i d \Delta_{n}}(\omega)\right)_{t \geq 0},\left(W_{2 i d \Delta_{n}+t}^{\prime}(\omega)-\right.\right.$ $\left.\left.W_{2 i d \Delta_{n}}^{\prime}(\omega)\right)_{t \geq 0}\right)$ similar to the functions occurring in Lemma 6.4, and since in the definition of $\xi_{i}^{n}$ one takes squared determinants, all these functions $\Phi$ are globally even in the sense that $\Phi(\omega, x, y)=\Phi(\omega,-x,-y)$ for any two functions $x, y$ on $\mathbb{R}_{+}$. So, on the one hand, after multiplying the function $\Phi$ corresponding to $\xi_{i}^{n}$ by $x^{m}\left(2 d \Delta_{n}\right)$ or $y^{m}\left(2 d \Delta_{n}\right)$, one gets an odd function, and (6.26) when $N=W^{m}$ or $N=W^{\prime m}$ follows. On the other hand, by the representation theorem one can write $\xi_{i}^{n}$ as the sum of two integrals over $\left(2 i d \Delta_{n}, 2(i+1) d \Delta_{n}\right]$ with respect to $W$ and $W^{\prime}$, for suitable predictable integrands; thus when $N$ is orthogonal to $W$ and $W^{\prime}$, the increment $N_{2(i+1) d \Delta_{n}}-N_{2 i d \Delta_{n}}$ has $\mathcal{F}_{i}^{n}$-conditional correlation 0 with both those integrals, thus yielding (6.26) again.

Therefore, the proof of Theorem 3.2 is complete.
(3) A simple calculation shows that

$$
\widehat{R}(n, T)-r=\frac{1}{\log 2} \log \frac{1+\sqrt{\Delta_{n}} U(r)_{T}^{n, 1} / S(r)_{T}}{1+\sqrt{\Delta_{n}} U(r)_{T}^{n, 2} / S(r)_{T}} \quad \text { if } S(r)_{T}>0
$$

hence on the set $\Omega_{T}^{r}$. Since the sequence $U(r)_{T}^{n}$ is tight, it follows from a Taylor expansion that

$$
\begin{equation*}
\frac{1}{\sqrt{\Delta_{n}}}(\widehat{R}(n, T)-r)-\frac{1}{S(r)_{T} \log 2}\left(U(r)_{T}^{n, 1}-U(r)_{T}^{n, 2}\right) \xrightarrow{\mathbb{P}} 0 \tag{6.27}
\end{equation*}
$$

on $\Omega_{T}^{r}$ again. Then Corollary 3.3 follows from Theorem 3.2, upon observing that the $\mathcal{F}$-conditional variance of $U(r)_{T}^{1}-U(r)_{T}^{2}$ is the numerator of the right-hand side of (3.13).
(4) Now we turn to the proof of (3.16), and by (6.19) it suffices to prove the convergence of $V(r)^{n, \kappa, \kappa^{\prime}}$. We suppose that $\kappa=\kappa^{\prime}$, the proof in the case $\kappa \neq \kappa^{\prime}$ being analogous. We set

$$
\eta_{i}^{n}=\gamma_{r}\left(A_{i}^{n, \kappa}, B_{i}^{n, \kappa}\right)^{4}, \quad \eta_{i}^{\prime n}=\mathbb{E}\left(\eta_{i}^{n} \mid \mathcal{F}_{i}^{n}\right), \quad \eta_{i}^{\prime \prime n}=\eta_{i}^{n}-\eta_{i}^{\prime n}
$$

As for (6.23), we deduce from (3.4) and (6.15) that

$$
\eta_{i}^{\prime n}=\Gamma_{r}^{\prime}\left(\underline{u}_{i}^{n}\right)+\Gamma_{r}\left(\underline{u}_{i}^{n}\right)^{2}, \quad \mathbb{E}\left(\left|\eta_{i}^{\prime \prime n}\right|^{2}\right) \leq K
$$

On the one hand, the same argument as for proving (6.25) shows that $4 d^{2} \Delta_{n} \times$ $\sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \eta_{i}^{\prime n}$ converges in the u.c.p. sense to the right-hand side of (3.16) (for $\kappa=\kappa^{\prime}$ ). On the other hand, since $\eta_{i}^{\prime \prime n}$ is a martingale increment relative to the filtration $\left(\mathcal{F}_{i}^{n}\right)_{i \geq 0}$, we deduce from Doob's inequality that $4 d^{2} \Delta_{n} \sum_{i=0}^{\left[t / 2 d \Delta_{n}\right]-1} \eta_{i}^{\prime \prime n} \xrightarrow{\text { u.c.p. }}$ 0 . We then deduce (3.16).
(5) Finally, for (3.17) it is enough to show the convergence in probability in restriction to each set $\Omega_{T}^{r}$, for $r=0, \ldots, d$. For this we use the following convergence properties, which readily follow from (3.8), (3.12), in restriction to the set $\Omega_{T}^{r}$ :

$$
\frac{1}{\Delta_{n}^{d-r}} S_{T}^{n, 1} \xrightarrow{\mathbb{P}} S(r)_{T}>0, \quad \widehat{R}(n, T) \xrightarrow{\mathbb{P}} r
$$

together with (3.16) applied at time $T$, which also holds on $\Omega_{T}^{r}$. Then (3.17) follows after a (slightly tedious) calculation, in view of the form (3.13) of $V(T)$ on $\Omega_{T}^{r}$ : the proof of Proposition 3.4 is complete.
6.5. Proofs of Theorems 4.1 and 4.5. We begin with a lemma. Its setting apparently extends the setting of the theorem to be proved, but this will be useful for the proof itself. The extension concerns the fact that we replace the nonrandom terminal time $T$ by a stopping time, still denoted by $T$, which is positive and bounded. In this case, the notation (2.3) still makes sense, as well as $A(p)_{T}^{n}$ and $a(n, T)$ and $\bar{V}_{T}^{n, \kappa \kappa^{\prime}}$, as given by (4.2) and (4.6).

Lemma 6.5. Assume Assumption (H) and $r_{t}=r$ for all $t \leq T$ with $T$ a positive finite stopping time and $r \in\{0, \ldots, d\}$. Then for all $p>0, \kappa, \kappa^{\prime} \in$ $\{1,2\}$ and $\Theta_{s}^{r, \kappa \kappa^{\prime}}$ as in (3.16) we have

$$
\begin{gather*}
A(p)_{T}^{n} \xrightarrow{\mathbb{P}} \operatorname{Tr}^{p},  \tag{6.28}\\
\frac{1}{\left(\kappa \kappa^{\prime} \Delta_{n}^{2}\right)^{d-r}} \bar{V}_{T}^{n, \kappa \kappa^{\prime}} \xrightarrow{\mathbb{P}} 2 d \int_{0}^{T}\left(\frac{1}{\Gamma_{r}\left(\sigma_{s}, \tilde{\sigma}, v_{s}, b_{s}\right)}-\frac{T}{S(r)_{T}}\right)^{2} \Theta_{s}^{r, \kappa \kappa^{\prime}} d s . \tag{6.29}
\end{gather*}
$$

Moreover, if $r \geq 1$, the following stable convergence in law holds, where $\mathcal{U}(r)$ is defined in Theorem 3.2:

$$
\begin{align*}
& \left(U(r)_{T}^{n}, \frac{1}{\sqrt{\Delta_{n}}}\left(A(p)_{T}^{n}-a(n, T) r^{p}\right)\right)  \tag{6.30}\\
& \quad \xrightarrow{\mathcal{L}-s}\left(\mathcal{U}(r)_{T}, \frac{p r^{p-1}}{\log 2} \int_{0}^{T} \frac{1}{\Gamma_{r}\left(\sigma_{s}, \tilde{\sigma}, v_{s}, b_{s}\right)}\left(d \mathcal{U}(r)_{s}^{1}-d \mathcal{U}(r)_{s}^{2}\right)\right)
\end{align*}
$$

Proof. (1) Let $\gamma_{t}=\Gamma_{r}\left(\sigma_{t}, \tilde{\sigma}, v_{t}, b_{t}\right)$, which is a continuous process, positive on $[0, T]$ by Lemma 3.1. Thus $T_{m}=m \wedge T \wedge \inf \left(t: \gamma_{t}<1 / m\right)$ satisfies $\mathbb{P}\left(T_{m}=T\right) \rightarrow 1$ as $m \rightarrow \infty$ and, if any one of the claimed convergence holds for each $T_{m}$ (instead of $T$ ), it also holds for $T$. In other words, we can assume $T \leq A$ and $1 / \gamma_{t} \leq A$ for some constant $A$ and all $t \in[0, T]$. Moreover, $\Gamma_{r}$ is a polynomial, so the process $\gamma_{t}$ is a continuous Itô semimartingale, and by localization again one can assume that for some other constant $A^{\prime}$,

$$
\begin{equation*}
\mathbb{E}\left(\left|\gamma_{t+s}-\gamma_{t}\right|^{2}\right) \leq A^{\prime} s \tag{6.31}
\end{equation*}
$$

The sequence $U(r)^{n}$ converges in law toward a continuous process, so the moduli of continuity $\rho(n, x)=\sup \left(\left\|U(r)_{t+s}^{n}-U(r)_{t}^{n}\right\|: t \leq A^{\prime},|s| \leq x\right)$ satisfy $\lim _{x \downarrow 0} \limsup _{n} \mathbb{P}(\rho(n, x)>1)=0$, and thus with the simplifying notation $w_{n}=2 d k_{n} \Delta_{n}$ we have

$$
\begin{align*}
& \mathbb{P}\left(\Omega_{n}\right) \rightarrow 1 \\
& \quad \text { where } \Omega_{n}=\left\{\left\|U(r)_{t+s}^{n}-U(r)_{t}^{n}\right\| \leq 1 \forall t \leq A^{\prime}, s \leq w_{n}\right\} . \tag{6.32}
\end{align*}
$$

(2) Observe that

$$
\widehat{R}_{i}^{n}=r+\frac{1}{\log 2} \log \frac{\zeta_{i}^{n}+\sqrt{\Delta_{n}} \eta_{i}^{n, 1}}{\zeta_{i}^{n}+\sqrt{\Delta_{n}} \eta_{i}^{n, 2}}
$$

where

$$
\zeta_{i}^{n}=S(r)_{2 i d \Delta_{n}+w_{n}}-S(r)_{2 i d \Delta_{n}}, \quad \eta_{i}^{n, k}=U(r)_{2 i d \Delta_{n}+w_{n}}^{n, k}-U(r)_{2 i d \Delta_{n}}^{n, k}
$$

Recalling $1 / \gamma_{t} \leq A$, and since $\zeta_{i}^{n}=w_{n}\left(\gamma_{2 i d \Delta_{n}}+\rho_{i}^{n}\right)$, where $\mathbb{E}\left(\left|\rho_{i}^{n}\right|^{2}\right) \leq A^{\prime \prime} w_{n}$ by (6.31), one has

$$
\begin{equation*}
0 \leq \frac{w_{n}}{\zeta_{i}^{n}} \leq A, \quad \mathbb{E}\left(\left|\frac{w_{n}}{\zeta_{i}^{n}}-\frac{1}{\gamma_{2 i d \Delta_{n}}}\right|^{2}\right) \leq A^{2} A^{\prime \prime} w_{n} \tag{6.33}
\end{equation*}
$$

Moreover, take $\alpha \in(0,1 / 2)$ such that $\frac{2|\log (1-\alpha)|}{\log 2} \leq 1 / 2$. For $n$ large enough, we have $A \sqrt{\Delta_{n}} / w_{n} \leq \alpha$ because $k_{n} \Delta_{n}^{3 / 4} \rightarrow \infty$. In this case, in restriction to the set $\Omega_{n}$, for all $i \leq\left[T / 2 d \Delta_{n}\right]-k_{n}-1$ we have with a constant $K$ (varying from place to place below):

$$
\begin{gather*}
\left|\frac{\sqrt{\Delta_{n}} \eta_{i}^{n, k}}{\zeta_{i}^{n}}\right| \leq \frac{A \sqrt{\Delta_{n}}}{w_{n}} \leq \frac{1}{2}, \\
\left|\widehat{R}_{i}^{n}-r-\frac{\sqrt{\Delta_{n}}}{\log 2} \frac{\eta_{i}^{n, 1}-\eta_{i}^{n, 2}}{\zeta_{i}^{n}}\right| \leq K \frac{\Delta_{n}}{w_{n}^{2}},  \tag{6.34}\\
r \geq 1 \Rightarrow\left|\frac{\widehat{R}_{i}^{n}}{r}-1\right| \leq \frac{K \sqrt{\Delta_{n}}}{w_{n}} \wedge \frac{1}{2}, \\
r=0 \Rightarrow\left|\widehat{R}_{i}^{n}\right| \leq \frac{K \sqrt{\Delta_{n}}}{w_{n}} \wedge \frac{1}{2} .
\end{gather*}
$$

(3) Recalling (4.2) and $T \leq A$, when $r=0$ the last estimate above yields

$$
\mathbb{E}\left(A(p)_{T}^{n} 1_{\Omega_{n}}\right) \leq K A \frac{\Delta_{n}^{p / 2}}{w_{n}^{p}}
$$

which goes to 0 because $k_{n} \sqrt{\Delta_{n}} \rightarrow \infty$. Thus in view of (6.32) one gets (6.28) when $r=0$.
(4) At this stage, we start proving (6.30), and thus assume $r \geq 1$. We observe that

$$
Y_{n}:=\frac{1}{\sqrt{\Delta_{n}}}\left(A(p)_{T}^{n}-a(n, T) r^{p}\right)=w_{n} \sum_{i=0}^{\left[T / w_{n}\right]-2} \xi_{i}^{n}
$$

where

$$
\xi_{i}^{n}=\frac{1}{\sqrt{\Delta_{n}}}\left(\left|\widehat{R}_{i k_{n}}^{n}\right|^{p} \wedge(d+1)^{p}-r^{p}\right) .
$$

Equation (6.34) implies that for $n$ large enough, $\left|\widehat{R}_{i}^{n}\right| \leq d+1$ (recall $r \leq d$ ), hence a Taylor expansion of the function $x \mapsto|r+x|^{p}-r^{p}$ implies, again for $n$ large enough:

$$
\left|\xi_{i}^{n}-\frac{p r^{p-1}}{\log 2} \frac{\eta_{i k_{n}}^{n, 1}-\eta_{i k_{n}}^{n, 2}}{\zeta_{i k_{n}}^{n}}\right| \leq K \frac{\sqrt{\Delta_{n}}}{w_{n}^{2}} \quad \text { on } \Omega_{n} \text { and for } i \leq\left[T / w_{n}\right]-2
$$

Upon using (6.33), and by the Cauchy-Schwarz inequality, it follows that

$$
\mathbb{E}\left(\left|\xi_{i}^{n}-\frac{p r^{p-1}}{\log 2} \frac{\eta_{i k_{n}}^{n, 1}-\eta_{i k_{n}}^{n, 2}}{w_{n} \gamma_{(i-1) w_{n}}}\right| 1_{\Omega_{n}}\right) \leq K \frac{\sqrt{\Delta_{n}}}{w_{n}^{2}}+K \sqrt{w_{n}}
$$

hence

$$
\mathbb{E}\left(\left|Y_{n}-Y_{n}^{\prime}\right| 1_{\Omega_{n}}\right) \rightarrow 0 \quad \text { where } Y_{n}^{\prime}=\frac{p r^{p-1}}{\log 2} \sum_{i=0}^{\left[T / w_{n}\right]-2} \frac{\eta_{i k_{n}}^{n, 1}-\eta_{i k_{n}}^{n, 2}}{\gamma_{(i-1) w_{n}}}
$$

because $k_{n} \Delta_{n}^{3 / 4} \rightarrow \infty$. Recall also (6.27). Then, by virtue of (6.32), the convergence (6.30) will follow from

$$
\begin{gather*}
\left(U(r)_{T}^{n}, Y_{n}^{\prime}\right) \xrightarrow{\mathcal{L}-s}\left(\mathcal{U}(r)_{T}, \mathcal{Y}\right),  \tag{6.35}\\
\mathcal{Y}=\frac{p r^{p-1}}{\log 2} \int_{0}^{T} \frac{1}{\Gamma_{r}\left(\sigma_{s}, \tilde{\sigma}, v_{s}, b_{s}\right)}\left(d \mathcal{U}(r)_{s}^{1}-d \mathcal{U}(r)_{s}^{2}\right) .
\end{gather*}
$$

(5) By Theorem VI.6.15 of [8] it follows from (6.23) and (6.25) that not only does the sequence of processes $U(r)^{n}$ converge in law, but it also enjoys the socalled P-UT property (predictable uniform tightness). By a trivial extension of Theorem VI.6.22 in [8], this implies that if a sequence $H^{n}$ of adapted càdlàg
two-dimensional processes on $\Omega$ is such that the pair $\left(U(r)^{n}, H^{n}\right) \xrightarrow{\mathcal{L}-s}(\mathcal{U}(r), H)$ (functional convergence for the Skorokhod topology), the bi-dimensional processes $\left(U(r)^{n}, \int_{0}^{t} H_{s-}^{n} d U(r)_{s}^{n}\right)$ converge stably in law to $\left(\mathcal{U}(r), \int_{0}^{t} H_{s-} d \mathcal{U}(r)_{s}\right)$, and since $\mathcal{U}(r)$ is continuous and $T$ is $\mathcal{F}$-measurable, this in turn implies the stable convergence of the variables $\left(U(r)_{T}^{n}, \int_{0}^{T} H_{s-}^{n} d U(r)_{s}^{n}\right) \xrightarrow{\mathcal{L}-s}\left(\mathcal{U}(r)_{T}\right.$, $\left.\int_{0}^{T} H_{s-} d \mathcal{U}(r)_{s}\right)$. At this point, (6.35) follows by taking the processes $H^{n}$ and $H$ with components

$$
\begin{aligned}
H_{t}^{n, 1} & =-H_{t}^{n, 2}=\frac{p r^{p-1}}{\gamma_{(i-1) w_{n}} \log 2} \quad \text { if } t \in\left((i-1) w_{n},\left(i w_{n}\right) \wedge T\right] \\
H_{t}^{1} & =-H_{t}^{2}=\frac{p r^{p-1}}{\gamma_{t} \log 2} 1_{\{t \leq T\}}
\end{aligned}
$$

[Note that the joint stable convergence $\left(U(r)^{n}, H^{n}\right) \xrightarrow{\mathcal{L}-s}(\mathcal{U}(r), H)$ holds because $1 / \gamma_{t}$ is continuous.] This ends the proof of (6.30).
(6) Since (6.30) implies (6.28) when $r \geq 1$, we are left to prove (6.29). We fix $\kappa, \kappa^{\prime}$. Our first observation is that, since $U(r)_{T}^{n} \xrightarrow{\mathcal{L}-s} \mathcal{U}(r)_{T}$ follows from (3.8) as seen before, the proof of (3.12) carries over to the case $T$ is a stopping time. Therefore $(\widehat{R}(n, T)-r) \log \Delta_{n} \xrightarrow{\mathbb{P}} 0$ because here $\Omega_{T}^{r}=\Omega$, and thus $\Delta_{n}^{\widehat{R}(n, T)-r} \xrightarrow{\mathbb{P}} 1$. It follows that (6.29) amounts to proving the same result for the variable $\widetilde{V}_{T}^{n, \kappa \kappa^{\prime}}$ which is the same as $\bar{V}_{T}^{n, \kappa \kappa^{\prime}}$ except that in front of the sum we substitute $\Delta_{n}^{1+2 d-2 \widehat{R}(n, T)}$ with $\Delta_{n}^{1+2 d-2 r}$.

With $\theta_{i}^{n}$ being as for Proposition 3.4, the $i$ th summand in the right-hand side of (3.15), we have

$$
\begin{align*}
& \widetilde{V}_{T}^{n, \kappa \kappa^{\prime}}=\sum_{j=0}^{2}\left(\frac{T \Delta_{n}^{d-r}}{S_{T}^{n, 1}}\right)^{j} B(j)_{T}^{n} \\
& \quad \text { where } B(j)_{T}^{n}=4 d^{2} \Delta_{n} \sum_{i=0}^{\left[T / 2 d \Delta_{n}\right]-k_{n}-1} v(j)_{i}^{n} \theta_{i}^{n}  \tag{6.36}\\
& v(0)_{i}^{n}=\left(\frac{w_{n}}{\zeta_{i}^{n}+\sqrt{\Delta_{n}} \eta_{i}^{n, 1}}\right)^{2}, \\
& v(1)_{i}^{n}=
\end{align*}
$$

Combining (6.33) and (6.34), we obtain for $i \leq\left[T / 2 d \Delta_{n}\right]-k_{n}-1$ and all $n$ large enough:

$$
\mathbb{E}\left(\left|v(0)_{i}^{n}-\frac{1}{\left(\gamma_{2 i d \Delta_{n}}\right)^{2}}\right|^{2} 1_{\Omega_{n}}\right)+\mathbb{E}\left(\left|v(1)_{i}^{n}+\frac{2}{\gamma_{2 i d \Delta_{n}}}\right|^{2} 1_{\Omega_{n}}\right) \leq K\left(\frac{\Delta_{n}}{w_{n}^{2}}+w_{n}\right) .
$$

Since by localization we may assume that the processes $\sigma_{t}, v_{t}, b_{t}$ are bounded, we may also assume $\theta_{i}^{n} \leq K$, and upon using (6.32) once more, we then deduce that $B(j)_{T}^{n}$ as the same asymptotic behavior as $B^{\prime}(j)_{T}^{n}$ which is given by the same formula, with $v(j)_{i}^{n}$ substituted with the following variables $v^{\prime}(j)_{i}^{n}$ :

$$
v^{\prime}(0)_{i}^{n}=\frac{1}{\left(\gamma_{2 i d \Delta_{n}}\right)^{2}}, \quad v^{\prime}(1)_{i}^{n}=-2 \frac{1}{\gamma_{2 i d \Delta_{n}}}, \quad v^{\prime}(2)_{i}^{n}=1 .
$$

This allows us to get, with $\Upsilon(0)=\frac{1}{\gamma^{2}}, \Upsilon(1)=\frac{1}{\gamma}$ and $\Upsilon(2)=1$,

$$
\begin{equation*}
\frac{1}{\left(\kappa \kappa^{\prime} \Delta_{n}^{2}\right)^{d-r}} B(j)_{T}^{n} \xrightarrow{\mathbb{P}} 2 d \int_{0}^{T} \Upsilon(j)_{s} \Theta_{s}^{r, \kappa \kappa^{\prime}} d s \tag{6.37}
\end{equation*}
$$

[indeed, the case $j=2$ is (3.16), and the other two cases follow from a standard argument, similar to step 5 above, but simpler because the integrand is a càdlàg bounded process not depending on $n$, and $\left.\theta_{i}^{n} \geq 0\right]$. Using further $T \Delta_{n}^{d-r} / S_{T}^{n, 1} \xrightarrow{\mathbb{P}}$ $T / S(r)_{T}$ and recalling (6.36), and upon expanding the square in the right-hand side of (6.29), we obtain this convergence and the lemma is proved.

PROOF OF (a) OF THEOREM 4.1. Since $\widehat{R}(n, T) \xrightarrow{\mathbb{P}} R_{T}$ by (3.12), it suffices to prove that

$$
\begin{equation*}
A(p)_{T}^{n} \xrightarrow{\mathbb{P}} A(p)_{T}:=\int_{0}^{T}\left(r_{s}\right)^{p} d s . \tag{6.38}
\end{equation*}
$$

The assumption implies the existence of a sequence of stopping times $\tau_{j}$ increasing to infinity, such that $\tau_{0}=0$ and $\tau_{j}<\tau_{j+1}$ if $\tau_{j}<\infty$, and such that the process $r_{t}$ takes a constant (random) value $\rho(j)$ on the time interval $J_{j}=\left(\tau_{j-1}, \tau_{j}\right)$, with $\rho(j) \neq \rho(j+1)$ if $0<\tau_{j}<\infty$. In view of the discussion preceding (2.3), the values $r_{\tau_{j}}$ is necessarily smaller than or equal to $\rho(j) \wedge \rho(j+1)$, but is irrelevant to our discussion. We also denote by $N_{T}$ the biggest $j$ such that $\tau_{j} \leq T$.

With an empty sum being set to 0 , we have

$$
A(p)_{T}^{n}=\sum_{j=1}^{N_{T}} Y(j)_{n}+Z_{n}, \quad Y(j)_{n}=w_{n} \sum_{i=\left[\tau_{j-1} / w_{n}\right]+1}^{\left[\left(\tau_{j} \wedge T\right) / w_{n}\right]-2}\left|\widehat{R}_{i}^{n}\right|^{p}
$$

and where $Z_{n}$ is the sum of at most $3 N_{T}$ terms of the form $w_{n}\left(\left|\widehat{R}_{i}^{n}\right|^{p} \wedge(d+1)^{p}\right)$. Since $N_{T}$ is finite and $w_{n} \rightarrow 0$, we have $Z_{t}^{n} \rightarrow 0$ (pointwise), and it suffices to show that for each $j \geq 1$ we have

$$
\begin{equation*}
Y(j)_{n} \xrightarrow{\mathbb{P}} Y(j):=\left(\left(T \wedge \tau_{j}\right)-\left(T \wedge \tau_{j-1}\right)\right) \rho(j)^{p} . \tag{6.39}
\end{equation*}
$$

We then fix $j$. The variable $Y(j)_{n}$ is the process $A(p)^{n}$ evaluated at time $T_{j}=$ $T \wedge \tau_{j}-\tau_{j-1}$ relative to the underlying process $X(j)_{t}=X_{\tau_{j-1}+t}$, up to at most two border terms. We thus might be tempted to apply (6.28) right away, and indeed
$X(j)$ satisfies Assumption $(\mathrm{H})$ for the filtration $\mathcal{F}(j)_{t}=\mathcal{F}_{\tau_{j-1}+t}$, relative to which $T_{j}$ is a positive bounded stopping time. There are, however, a few problems to overcome:

1. The rank $r_{t}(X(j))$ associated with $X(j)$ is equal to $\rho(j)$ for all $t \in\left(0, T_{j}\right)$, but not necessarily for $t=0$, nor for $t=T_{j}$;
2. This rank $\rho(j)$ is random, albeit $\mathcal{F}(j)_{0}$-measurable.

To solve these problems, we fix $\varepsilon>0$ and consider the process $X(j, \varepsilon)_{t}=$ $X_{\tau_{j-1}+\varepsilon+t}$, satisfying Assumption (H) for the filtration $\mathcal{F}(j, \varepsilon)_{t}=\mathcal{F}_{\tau_{j-1}+\varepsilon+t}$, and the $\mathcal{F}(j, \varepsilon)_{t}$-stopping time $T(j, \varepsilon)=T \wedge \tau_{j}-T \wedge \tau_{j-1}-2 \varepsilon$. The associated rank is thus $\rho(j)$ for all $t \in[0, T(j, \varepsilon)]$, and we will show that if $A(p, j, \varepsilon)$ is associated with $X(j, \varepsilon)$ by (4.2), we have

$$
\begin{equation*}
A(p, j, \varepsilon)_{T(j, \varepsilon)}^{n} \xrightarrow{\mathbb{P}} T(j, \varepsilon) \rho(j)^{p}=\left(T \wedge \tau_{j}-T \wedge \tau_{j-1}-2 \varepsilon\right) \rho(j)^{p} \tag{6.40}
\end{equation*}
$$

Indeed, it suffices to prove this in restriction to each set $\Omega_{r}^{\prime}=\{\rho(j)=r\}$ satisfying $\mathbb{P}\left(\Omega_{r}^{\prime}\right)>0$. If $\mathbb{P}_{r}$ denotes the (usual) conditional probability $\mathbb{P}\left(\cdot \mid \Omega_{r}^{\prime}\right)$, the process $X(j, \varepsilon)$, on the space $\left(\Omega_{r}^{\prime}, \mathcal{F} \cap \Omega_{r}^{\prime},\left(\mathcal{F}(j, \varepsilon)_{t} \cap \Omega_{r}^{\prime}\right), \mathbb{P}_{r}\right)$, still satisfies Assumption (H) and the associated rank is now $r$ on the time interval $[0, T(j, \varepsilon)]$. Then Lemma 6.5 yields the convergence (6.40) under $\mathbb{P}_{r}$, hence under $\mathbb{P}$ in restriction to each $\Omega_{r}^{\prime}$, hence under $\mathbb{P}$ on $\Omega$ itself.

Finally, the difference $Y(j)_{n}-A(p, j, \varepsilon)_{T(j, \varepsilon)}^{n}$ is a sum of at most $2\left[\varepsilon / w_{n}\right]$ terms, each one smaller than $w_{n}(d+1)^{p}$, so this difference is smaller than $K \varepsilon$, as is the difference between the two right sides of (6.39) and (6.40). Hence (6.39) follows from (6.40), by taking first $n \rightarrow \infty$ and then $\varepsilon \rightarrow 0$. This completes the proof.

Proof of (b) OF Theorem 4.1. Exactly as in the previous proof, it is enough to prove the result when $r_{t}=r \geq 1$ identically, for some nonrandom $r \in\{1, \ldots, d\}$. By a standard localization procedure, we can assume that $\Gamma_{r}\left(\sigma_{t}, \tilde{\sigma}, v_{t}, b_{t}\right)$, which is positive everywhere, is bounded from below by a constant $1 / A$ with $A>0$, so the assumptions of Lemma 6.5 are satisfied. Therefore, (6.27) and (6.30) yield that, with $Y_{n}$ and $\mathcal{Y}$ as in the proof of Lemma 6.5 and with $Z_{n}=\frac{1}{\sqrt{\Delta_{n}}}(\widehat{R}(n, T)-r)$ and $\mathcal{Z}=\frac{1}{S(r)_{T} \log 2}\left(\mathcal{U}(r)_{T}^{1}-\mathcal{U}(r)_{T}^{2}\right)$, we have

$$
\begin{equation*}
\left(Y_{n}, Z_{n}\right) \xrightarrow{\mathcal{L}-s}(\mathcal{Y}, \mathcal{Z}) \tag{6.41}
\end{equation*}
$$

Then we obtain

$$
\frac{1}{\sqrt{\Delta_{n}}} B(n, p, T)=Y_{n}+\frac{a(n, T)}{\sqrt{\Delta_{n}}}\left(r^{p}-\left|r+\sqrt{\Delta_{n}} Z_{n}\right|^{p}\right) .
$$

On the one hand, $a(n, T) \rightarrow T$. On the other hand, since $Z_{n}$ converges in law and $r \geq 1$, we have by the mean value theorem

$$
\frac{1}{\sqrt{\Delta_{n}}}\left(r^{p}-\left|r+\sqrt{\Delta_{n}} Z_{n}\right|^{p}\right)+p r^{p-1} Z_{n} \xrightarrow{\mathbb{P}} 0 .
$$

Hence, (6.41) yields

$$
\frac{1}{\sqrt{\Delta_{n}}} B(n, p, T) \xrightarrow{\mathcal{L}-\mathrm{s}} \mathcal{B}(p, T)=\mathcal{Y}-T p r^{p-1} \mathcal{Z} .
$$

The pair $(\mathcal{Y}, \mathcal{Z})$ being $\mathcal{F}$-conditionally centered Gaussian, the same is true of $\mathcal{B}(p, T)$, and the form (4.5) of its conditional variance is easily checked by virtue of (3.9).

Proof of Theorem 4.5. It is easy to construct a process $X^{\prime}$ which satisfies the assumptions of (a) of Theorem 4.1 and such that $X_{t}^{\prime}=X_{t}$ for all $t \leq T$ on the set $\Omega_{T}^{\neq}$. Then on this set $B(n, p, T)$ is the same, when constructed upon $X$ or upon $X^{\prime}$, and thus it converges in probability on this set to a strictly negative variable. On the other hand, $Z(n, p, T)$ is $B(n, p, T)$ divided by a quantity which by construction is smaller than $\Delta_{n}^{1 / 4}$. Then the convergence $Z(n, p, T) \rightarrow-\infty$ on $\Omega_{T}^{\neq}$is clear.

It suffices to prove (4.7) on the set $\Omega_{T}^{\overline{=}} \cap \Omega_{T}^{r}$ for any $r \in\{0, \ldots, d\}$ such that $\mathbb{P}\left(\Omega_{T}^{r}\right)>0$. For this, we can argue under the conditional probability $\mathbb{P}_{r}=\mathbb{P}(\cdot \mid$ $\left\{r_{0}=r\right\}$ ), or equivalently suppose that we have in fact $r_{0}=r$. As above, one can construct a process $X^{\prime}$ which satisfies the assumptions of Lemma 6.5 for some stopping time $T^{\prime}$ which satisfies $T^{\prime} \geq T$ on the set $\Omega_{T}^{=}$, and we can apply (6.29) to $X^{\prime}$ and the stopping time $T^{\prime} \wedge T$. This gives us (6.29) for $X$, in restriction to the set $\Omega_{T}^{=}$.

At this point, (4.7) follows from (6.29) by exactly the same calculations as (3.17) follows from (3.16).

Finally, since $\bar{V}(n, p, T) \xrightarrow{\mathbb{P}} \bar{V}(p, T)$ on $\Omega_{\bar{T}}^{\overline{=}}$, we have $Z(n, p, T)=B(n, p$, $T) / \sqrt{\Delta_{n} B V(n, p, T)}$ on a set $\Omega_{n}^{\prime \prime}$ whose probability goes to 1 . The first part of (4.9) than follows from (4.4) and (4.7) by delta method for stable convergence.
6.6. Proof of Corollary 3.6. The same stopping argument as in step 2 of the previous proof allows us to show that, without assumptions on the rank process $r_{t}$, the stable convergence in law (3.8) holds in restriction to the set $\Omega_{T}^{r}$, as soon as we restrict our attention to the time interval $[0, T]$.

At this stage, the first claim of Corollary 3.6 follows from (3.12), an application of the delta method, (3.17) and classical properties of stable convergence in law. Equation (3.14) is in turn a trivial consequence of (3.12) and the definition of $\widehat{R}^{\mathrm{int}}(n, T)$.
6.7. Proof of Propositions 3.7, 3.8 and 4.6. Equation (3.20) is an obvious consequence of the stable convergence (3.18). For the alternative-consistency, it suffices to prove that for any $r^{\prime} \neq r$ we have

$$
\begin{equation*}
\mathbb{P}\left(\mathcal{C}(\alpha)_{T}^{n,=r} \cap \Omega_{T}^{r^{\prime}}\right) \rightarrow \mathbb{P}\left(\Omega_{T}^{r^{\prime}}\right) \tag{6.42}
\end{equation*}
$$

On the set $\Omega_{T}^{r^{\prime}}$ we have $S(n, T) \xrightarrow{\mathbb{P}} 2^{d-r^{\prime}}$, and by (3.17) the variables $V(n, T)$ converge in probability to a limit which is $[0, \infty)$-valued (actually, it is a.s. positive, but we do not use this fact here), so that $\Delta_{n} V(n, T) \xrightarrow{\mathbb{P}} 0$. Since $r^{\prime} \neq r$, (6.42) readily follows from the definition of $\mathcal{C}(\alpha)_{T}^{n,=r}$.

Propositions 3.8 and 4.6 are proved analogously, the alternative-consistency in the latter case following from the second part of (4.9).

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## REFERENCES

[1] Barndorff-Nielsen, O. E., Graversen, S. E., Jacod, J., Podolskij, M. and ShepHARD, N. (2006). A central limit theorem for realised power and bipower variations of continuous semimartingales. In From Stochastic Calculus to Mathematical Finance. Festschrift in Honour of A. N. Shiryaev (Yu. Kabanov, R. Liptser and J. Stoyanov, eds.) 33-68. Springer, Berlin. MR2233534
[2] Cont, R. and Mancini, C. (2011). Nonparametric tests for pathwise properties of semimartingales. Bernoulli 17 781-813. MR2787615
[3] Cragg, J. G. and Donald, S. G. (1996). On the asymptotic properties of LDU-based tests of the rank of a matrix. J. Amer. Statist. Assoc. 91 1301-1309. MR1424626
[4] Delbaen, F. and Schachermayer, W. (1994). A general version of the fundamental theorem of asset pricing. Math. Ann. 300 463-520. MR1304434
[5] Ditlevsen, S. and SøRENSEN, M. (2004). Inference for observations of integrated diffusion processes. Scand. J. Stat. 31 417-429. MR2087834
[6] Gloter, A. (2006). Parameter estimation for a discretely observed integrated diffusion process. Scand. J. Stat. 33 83-104. MR2255111
[7] JACOD, J. (2008). Asymptotic properties of realized power variations and related functionals of semimartingales. Stochastic Process. Appl. 118 517-559. MR2394762
[8] Jacod, J. and Shiryaev, A. N. (2002). Limit Theorems for Stochastic Processes, 2nd ed. Springer, Berlin. MR0959133
[9] Jacod, J., Lejay, A. and Talay, D. (2008). Estimation of the Brownian dimension of a continuous Itô process. Bernoulli 14 469-498. MR2544098
[10] Jacod, J. and Protter, P. (2012). Discretization of Processes. Springer, Heidelberg. MR2859096
[11] Kritchman, S. and Nadler, B. (2008). Determining the number of components in a factor model from limited noisy data. Chemometr. Intell. Lab. Syst. 94 19-32.
[12] Podolskij, M. and Rosenbaum, M. (2012). Testing the local volatility assumption: A statistical approach. Ann. Finance 8 31-48. MR2881178
[13] Robin, J.-M. and Smith, R. J. (2000). Tests of rank. Econometric Theory 16 151-175. MR1763430
[14] Vetter, M. (2012). Estimation of correlation for continuous semimartingales. Scand. J. Stat. 39 757-771. MR3000847

Institut de Mathématiques de Jussieu
CNRS-UMR 7586
Université Pierre et Marie Curie
4 Place Jussieu, Paris 75005
France
E-MAIL: jean.jacod@upmc.fr

Department of Mathematics
Heidelberg University
INF 294, HEIDELBERG 69120
GERMANY
E-MAIL: m.podolskij@uni-heidelberg.de


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