# ASYMPTOTIC PROPERTIES OF THE MAXIMUM LIKELIHOOD ESTIMATOR FOR STOCHASTIC PARABOLIC EQUATIONS WITH ADDITIVE FRACTIONAL BROWNIAN MOTION 

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A parameter estimation problem is considered for a diagonalizable stochastic evolution equation using a finite number of the Fourier coefficients of the solution. The equation is driven by additive noise that is white in space and fractional in time with the Hurst parameter $H \geq 1 / 2$. The objective is to study asymptotic properties of the maximum likelihood estimator as the number of the Fourier coefficients increases. A necessary and sufficient condition for consistency and asymptotic normality is presented in terms of the eigenvalues of the operators in the equation.

Keywords: Asymptotic normality; ergodicity; parameter estimation; stochastic evolution equations.

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## 1. Introduction

In the classical statistical estimation problem, the starting point is a family $\mathbf{P}^{\theta}$ of probability measures depending on the parameter $\theta$ in some subset $\Theta$ of a finitedimensional Euclidean space. Each $\mathbf{P}^{\theta}$ is the distribution of a random element. It is assumed that a realization of one random element corresponding to one value $\theta=\theta_{0}$ of the parameter is observed, and the objective is to estimate the values of this parameter from the observations.

The intuition is to select the value $\theta$ corresponding to the random element that is most likely to produce the observations. A rigorous mathematical implementation of this idea leads to the notion of the regular statistical model [4]: the statistical model (or estimation problem) $\mathbf{P}^{\theta}, \theta \in \Theta$, is called regular, if there exists a probability measure $\mathbf{Q}$ such that all the measures $\mathbf{P}^{\theta}$ are absolutely continuous with respect to $\mathbf{Q}$ and the density $d \mathbf{P}^{\theta} / d \mathbf{Q}$, called the likelihood ratio, has a certain regularity.

In regular models, the estimator $\widehat{\theta}$ of the unknown parameter is constructed by maximizing the likelihood ratio and is called the maximum likelihood estimator (MLE). Since, as a rule, $\widehat{\theta} \neq \theta_{0}$, the consistency of the estimator is studied, that is, the convergence of $\widehat{\theta}$ to $\theta_{0}$ as more and more information becomes available. In all known regular statistical problems, the amount of information can be increased in one of two ways: (a) increasing the sample size, for example, the observation time interval (large sample asymptotic); (b) reducing the amplitude of noise (small noise asymptotic).

In finite-dimensional models, the only way to increase the sample size is to increase the observation time. In infinite-dimensional models, in particular, those provided by stochastic partial differential equations (SPDEs), another possibility is to increase the dimension of the spatial projection of the observations. Thus, a consistent estimator can be possible on a finite time interval with fixed noise intensity. This possibility was first suggested by Huebner et al. [2] for parabolic equations driven by additive spacetime white noise, and was further investigated by Huebner and Rozovskii [3], where a necessary and sufficient condition for the existence of a consistent estimator was stated in terms of the orders of the operators in the equation.

The objective of this paper is to extend the model from [3] to parabolic equations in which the time component of the noise is fractional with the Hurst parameter $H \geq 1 / 2$. More specifically, we consider an abstract evolution equation

$$
\begin{equation*}
u(t)+\int_{0}^{t}\left(\mathcal{A}_{0}+\theta \mathcal{A}_{1}\right) u(s) d s=W^{H}(t) \tag{1.1}
\end{equation*}
$$

where $\mathcal{A}_{0}$ and $\mathcal{A}_{1}$ are known linear operators and $\theta \in \Theta \subseteq \mathbb{R}$ is the unknown parameter; the zero initial condition is taken to simplify the presentation. The noise $W^{H}(t)$ is the cylindrical fractional Brownian motion on a Hilbert space $\mathbf{H}$ and can be written as

$$
\begin{equation*}
W^{H}(t)=\sum_{j=1}^{\infty} w_{j}^{H}(t) h_{j}, \tag{1.2}
\end{equation*}
$$

where $\left\{w_{j}^{H}, j \geq 1\right\}$ are independent fractional Brownian motions with the same Hurst parameter $H \geq 1 / 2$ and $\left\{h_{j}, j \geq 1\right\}$ is an orthonormal basis in a Hilbert space $\mathbf{H} ; H=1 / 2$ corresponds to the cylindrical Brownian motion. It can be shown (see Proposition 2.1 below) that (1.2) defines a continuous Gaussian process with values in some larger Hilbert space X. Existence and uniqueness of the solution under certain regularity conditions on $\mathcal{A}_{0}$ and $\mathcal{A}_{1}$ for such
equations are well known for all $H \in(0,1)$; see, for example, Tindel et al. [17, Theorem 1].

The main additional assumption about (1.1), both in [3] and in this paper, is that the equation is diagonalizable: $\left\{h_{j}, j \geq 1\right\}$ from (1.2) is a common system of eigenfunction of the operators $\mathcal{A}_{0}$ and $\mathcal{A}_{1}$ :

$$
\begin{equation*}
\mathcal{A}_{0} h_{j}=\rho_{j} h_{j}, \quad \mathcal{A}_{1} h_{j}=\nu_{j} h_{j} . \tag{1.3}
\end{equation*}
$$

Under certain conditions on the numbers $\rho_{j}, \nu_{j}$, the solution of (1.1) is a convergent Fourier series $u(t)=\sum_{j \geq 1} u_{j}(t) h_{j}$, and each $u_{j}(t)$ is a fractional OrnsteinUhlenbeck (OU) process. An $N$-dimensional projection of the solution is then an N dimensional fractional OU process with independent components. A Girsanov-type formula (for example, from Kleptsyna et al. [8, Theorem 3]) leads to a maximum likelihood estimator $\hat{\theta}_{N}$ of $\theta$ based on the first $N$ Fourier coefficients $u_{1}, \ldots, u_{N}$ of the solution of (1.1). An explicit expression for this estimator exists but requires a number of additional notations; see formula (3.8) below. A particular case of Eq. (1.1) was studied by Prakasa Rao [14].

If the solution of (1.1) is completely observable at every $t \in[0, T]$, then each Fourier coefficient $u_{k}$, being a linear functional of the solution, can be computed. For example, if $u$ is regular enough so that $u(t) \in L_{2}(\Omega ; H)$ for every $t$, then $u_{k}(t)=\left(u(t), h_{k}\right)_{H}$.

The following is the main result of the paper.

Theorem 1.1. Define $\mu_{j}=\theta \nu_{j}+\rho_{j}$. Then the maximum likelihood estimator $\hat{\theta}_{N}$ of $\theta$ is strongly consistent and asymptotically normal, as $N \rightarrow \infty$, if and only if the series $\sum_{j} \nu_{j}^{2} \mu_{j}^{-1}$ diverges; the rate of convergence of the estimator is given by the square root of the partial sums of this series: as $N \rightarrow \infty$, the sequence $\left(\sum_{j \leq N} \nu_{j}^{2} \mu_{j}^{-1}\right)^{1 / 2}\left(\hat{\theta}_{N}-\theta\right)$ converges in distribution to a Gaussian random variable.

If the operators $\mathcal{A}_{0}$ and $\mathcal{A}_{1}$ are elliptic of orders $m_{0}$ and $m_{1}$ in $L_{2}(M)$, where $M$ is a $d$-dimensional manifold, and $2 m=\max \left(m_{0}, m_{1}\right)$, then the asymptotic of the eigenvalues of elliptic operators implies that the divergence of $\sum_{j} \nu_{j}^{2} \mu_{j}^{-1}$ is equivalent to $m_{1} \geq m-(d / 2)$; in the case $H=1 / 2$ this condition on the orders was derived in [3]. Thus, beside extending the results of [3] to fractional-in-time noise, we also generalize the necessary and sufficient condition for consistency of the estimator.

While parameter estimation for the finite-dimensional fractional OU and similar processes has been recently investigated by Tudor and Viens [18] for all $H \in(0,1)$, our analysis in infinite dimensions requires more delicate results: an explicit expression for the Laplace transform of a certain functional of the fractional OU process, as obtained by Kleptsyna and Le Breton [7], and for now this expression exists only for $H \geq 1 / 2$.

## 2. Stochastic Parabolic Equations with Additive FBM

In this section we introduce a diagonalizable stochastic parabolic equation driven by a cylindrical fractional Brownian motion and study the main properties of the solution.

Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space and let $\mathbf{H}$ be a separable Hilbert space with an orthonormal basis $\left\{h_{j}, j \geq 1\right\}$.

Definition 2.1. (a) A fractional Brownian motion with a Hurst parameter $H \in$ $(0,1)$ is a Gaussian process $w^{H}$ with zero mean and covariance

$$
\mathbb{E} w^{H}(t) w^{H}(s)=\frac{1}{2}\left(t^{2 H}+s^{2 H}-|t-s|^{2 H}\right), \quad t, s \geq 0
$$

(b) The cylindrical fractional Brownian motion $W^{H}$ on a Hilbert space $\mathbf{H}$ is a collection of zero-mean Gaussian processes $W_{f}^{H}=W_{f}^{H}(t)$ indexed by the elements $f$ of the space $\mathbf{H}$ such that, for every $f, g \in \mathbf{H}$ and $t, s \geq 0$,

$$
\begin{equation*}
\mathbb{E}\left(W_{f}^{H}(t) W_{g}^{H}(s)\right)=\frac{(f, g)_{\mathbf{H}}}{2}\left(t^{2 H}+s^{2 H}-|t-s|^{2 H}\right) . \tag{2.1}
\end{equation*}
$$

Proposition 2.1. Let $\mathbf{H}$ be a separable Hilbert space with an orthonormal basis $\left\{h_{j}, j \geq 1\right\}$.
(a) If $W^{H}$ is a cylindrical Brownian motion on $\mathbf{H}$, then $\left\{W_{h_{j}}^{H}(t), j \geq 1\right\}$ is a collection of independent fractional Brownian motions with the same Hurst parameter $H$.
(b) Let $\mathbf{X}$ is a Hilbert space such that $\mathbf{H}$ is a dense subset of $\mathbf{X}$ and the inclusion operator $\mathfrak{j}: \mathbf{H} \rightarrow \mathbf{X}$ is Hilbert-Schmidt. If $\left\{w_{j}^{H}(t), j \geq 1\right\}$ is a collection of independent fractional Brownian motions with the same Hurst parameter $H$, then the process $W^{H}(t)=\sum_{j \geq 1} w_{j}^{H}(t) h_{j}$ is a continuous $\mathbf{X}$-valued Gaussian process and is a cylindrical fractional Brownian motion on $\mathbf{H}$.

Proof. (a) This follows directly from (2.1).
(b) Since the inclusion $\mathfrak{j}: \mathbf{H} \hookrightarrow \mathbf{X}$ is Hilbert-Schmidt, we have

$$
\sum_{j \geq 1}\left\|h_{j}\right\|_{\mathbf{X}}^{2}=C_{\mathrm{j}}<\infty
$$

Therefore, since

$$
\mathbb{E} w_{j}^{H}(t) w_{m}^{H}(t)= \begin{cases}t^{2 H}, & \text { if } k=m \\ 0, & \text { if } k \neq m\end{cases}
$$

we find

$$
\mathbb{E}\left\|W^{H}(t)\right\|_{\mathbf{X}}^{2}=t^{2 H} C_{\mathbf{j}}<\infty
$$

Similarly,

$$
\mathbb{E}\left\|W^{H}(t)-W^{H}(s)\right\|_{\mathbf{X}}^{2}=|t-s|^{2 H} C_{\mathbf{j}},
$$

which, for $H>1 / 2$, implies continuity by the Kolmogorov criterion. For $H \leq 1 / 2$, we consider $\mathbb{E}\left\|W^{H}(t)-W^{H}(s)\right\|_{\mathbf{X}}^{2 n}$ for sufficiently large integer $n$. Next, setting

$$
W_{f}^{H}(t)=\sum_{j \geq 1}\left(f, h_{j}\right)_{\mathbf{H}}^{2}
$$

we see that $W_{f}^{H}$ is a zero-mean Gaussian process and (2.1) holds.
Consider the following equation:

$$
\begin{equation*}
d u(t)+\left(\mathcal{A}_{0}+\theta \mathcal{A}_{1}\right) u(t) d t=d W^{H}(t), \quad 0<t \leq T \tag{2.2}
\end{equation*}
$$

where $\mathcal{A}_{0}, \mathcal{A}_{1}$ are linear operators, $\theta$ is a scalar parameter belonging to an open set $\Theta \subset \mathbb{R}$, and, for simplicity, $u(0)=0$.

Definition 2.2. Equation (2.2) is called diagonalizable if the operators $\mathcal{A}_{0}$ and $\mathcal{A}_{1}$ have point spectrum and a common system of eigenfunctions $\left\{h_{j}, j \geq 1\right\}$.

Denote by $\rho_{j}, \nu_{j}$, and $\mu_{j}(\theta)$ the eigenvalues of the operators $\mathcal{A}_{0}, \mathcal{A}_{1}$, and $\mathcal{A}_{0}+\theta \mathcal{A}_{1}$ :

$$
\begin{equation*}
\mathcal{A}_{0} h_{j}=\rho_{j} h_{j}, \quad \mathcal{A}_{1} h_{j}=\nu_{j} h_{j}, \quad \mu_{j}(\theta)=\rho_{j}+\theta \nu_{j}, \quad j \geq 1 \tag{2.3}
\end{equation*}
$$

Definition 2.3. A diagonalizable equation (2.2) is called parabolic if there exist positive numbers $C^{*}, c_{1}, c_{2}$ such that, after possible re-arrangement, $\left\{\mu_{j}(\theta)+\right.$ $\left.C^{*}, j \geq 1\right\}$ is a positive, non-decreasing, and unbounded sequence for all $\theta \in \Theta$ and

$$
\begin{equation*}
c_{1} \leq \frac{\mu_{j}\left(\theta_{1}\right)+C^{*}}{\mu_{j}\left(\theta_{2}\right)+C^{*}} \leq c_{2} \tag{2.4}
\end{equation*}
$$

for all $\theta_{1}, \theta_{2} \in \Theta$.
In particular, if Eq. (2.2) is parabolic, then $\lim _{j \rightarrow \infty} \mu_{j}(\theta)=+\infty$ for all $\theta \in \Theta$, and there exists an index $J \geq 1$ such that, for all $\theta \in \Theta$ and $j \geq J$,

$$
\begin{equation*}
\mu_{j}(\theta)>0 \tag{2.5}
\end{equation*}
$$

Indeed, if such index does not exist, then there exists a sequence of $\theta_{j} \in \Theta$, such that $\mu_{j}\left(\theta_{j}\right) \leq 0$, and hence $\mu_{j}\left(\theta_{j}\right)+C^{*} \leq C^{*}$. By (2.4), for a fixed $\theta$, we have

$$
\lim _{j \rightarrow \infty} \frac{\mu_{j}\left(\theta_{j}\right)+C^{*}}{\mu_{j}(\theta)+C^{*}}=0
$$

that contradicts the assumption that $c_{1}>0$.
Example 2.1. Let $G$ be a smooth bounded domain in $\mathbb{R}^{d}$ or a smooth compact $d$ dimensional manifold with a smooth measure, $\mathbf{H}=L_{2}(G)$, and let $\boldsymbol{\Delta}$ be the Laplace operator on $G$ (with zero boundary conditions if $G$ is a domain). It is known (see, for example, Safarov and Vassiliev [15] or Shubin [16]) that $\boldsymbol{\Delta}$ has a complete orthonormal system of eigenfunctions in $\mathbf{H}$, and the corresponding eigenvalues $\lambda_{j}$
are negative, can be arranged in decreasing order, and there is a positive number $c_{\text {o }}$ such that

$$
\begin{equation*}
\lim _{j \rightarrow \infty} j^{-2 / d}\left|\lambda_{j}\right|=c_{\circ} \tag{2.6}
\end{equation*}
$$

The reader can verify that each of the following equations is diagonalizable and parabolic:

$$
\begin{gather*}
d u-\theta \boldsymbol{\Delta} u d t=d W^{H}, \theta \geq a>0  \tag{2.7}\\
d u-(\boldsymbol{\Delta} u+\theta u) d t=d W^{H}  \tag{2.8}\\
d u+\left(\boldsymbol{\Delta}^{2} u+\theta \boldsymbol{\Delta} u\right) d t=d W^{H} \tag{2.9}
\end{gather*}
$$

The following result both defines the solution of (2.2) and establishes existence and uniqueness of the solution.

Theorem 2.1. Assume that Eq. (2.2) is diagonalizable and parabolic, $W^{H}$ is a cylindrical fractional Brownian motion on a separable Hilbert space $\mathbf{H},\left\{h_{j}, j \geq 1\right\}$ is an orthonormal basis in $\mathbf{H}$, and $H \geq 1 / 2$.

If $\mathbf{X}$ is a Hilbert space such that $\mathbf{H}$ is a dense subset of $\mathbf{X}$ and the inclusion operator $\mathfrak{j}: \mathbf{H} \rightarrow \mathbf{X}$ is Hilbert-Schmidt, then the process $u$ defined by

$$
\begin{equation*}
u(t)=\sum_{j \geq 1} u_{j}(t) h_{j} \tag{2.10}
\end{equation*}
$$

where

$$
\begin{equation*}
u_{j}(t)=\int_{0}^{t} e^{-\mu_{j}(\theta)(t-s)} d W_{h_{j}}^{H}(s) \tag{2.11}
\end{equation*}
$$

is a Gaussian process with values in $\mathbf{X}$.
Proof. The argument is similar to the proof of Proposition 2.1(b). The properties of the fractional Brownian motion with $H>1 / 2$ imply

$$
\mathbb{E} u_{j}^{2}(t)=H(2 H-1) e^{-2 \mu_{j}(\theta) t} \int_{0}^{t} \int_{0}^{t} e^{\mu_{j}(\theta)\left(s_{1}+s_{2}\right)}\left|s_{1}-s_{2}\right|^{2 H-2} d s_{1} d s_{2}
$$

see, for example, Pipiras and Taqqu [13, formulas (4.1), (4.2)]. By direct computations,

$$
\begin{aligned}
\mathbb{E} u_{j}^{2}(t) & =\frac{H(2 H-1)}{\mu_{j}(\theta)}\left(\int_{0}^{t} s^{2 H-2} e^{-\mu_{j}(\theta) s} d s-e^{-2 \mu_{j}(\theta) t} \int_{0}^{t} s^{2 H-2} e^{\mu_{j}(\theta) s} d s\right) \\
& =\frac{H(2 H-1)}{\left|\mu_{j}(\theta)\right|^{2 H}}\left(\int_{0}^{\mu_{j}(\theta) t} s^{2 H-2} e^{-s} d s-e^{-\mu_{j}(\theta) t} \int_{0}^{t} s^{2 H-2} e^{-\mu_{j}(\theta)(t-s)} d s\right)
\end{aligned}
$$

As $\lim _{j \rightarrow \infty} \mu_{j}(\theta)=+\infty$, we find, for every $t>0$,

$$
\begin{aligned}
\lim _{j \rightarrow \infty} \int_{0}^{\mu_{j}(\theta) t} s^{2 H-2} e^{-s} d s=\int_{0}^{\infty} s^{2 H-2} e^{-s} d s=\Gamma(2 H-1), \\
\lim _{j \rightarrow \infty}\left|e^{-\mu_{j}(\theta) t} \int_{0}^{t} s^{2 H-2} e^{-\mu_{j}(\theta)(t-s)} d s\right| \leq \lim _{j \rightarrow \infty} \frac{e^{-\mu_{j}(\theta) t} t^{2 H-1}}{2 H-1}=0
\end{aligned}
$$

Hence,

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left|\mu_{j}(\theta)\right|^{2 H} \mathbb{E} u_{j}^{2}(t)=H(2 H-1) \Gamma(2 H-1) \tag{2.12}
\end{equation*}
$$

and the limit indeed does not depend on $t$ as long as $t>0$. It follows from (2.12) that $\sup _{j} \mathbb{E}\left|u_{j}(t)\right|^{2}<\infty$ and therefore

$$
\begin{equation*}
\mathbb{E}\|u(t)\|_{\mathbf{X}}^{2}=\sum_{j=1}^{\infty} \mathbb{E}\left|u_{j}(t)\right|^{2}\left\|h_{j}\right\|_{\mathbf{X}}^{2}<\infty \tag{2.13}
\end{equation*}
$$

If $H=1 / 2$, then $W_{h_{j}}^{H}$ is a standard Brownian motion, so that

$$
\mathbb{E} u_{j}^{2}=\int_{0}^{t} e^{-2 \mu_{j}(\theta)(t-s)} d s \quad \text { and } \quad \lim _{j \rightarrow \infty}\left|\mu_{j}(\theta)\right| \mathbb{E} u_{j}^{2}=\frac{1}{2}
$$

note that passing to the limit $H \rightarrow 1 / 2$ in (2.12) gives the same result.

Definition 2.4. The process $u$ constructed in Theorem 2.1 is called the solution of Eq. (2.2).

## 3. The Maximum Likelihood Estimator and its Properties

Consider the diagonalizable equation

$$
\begin{equation*}
d u(t)+\left(\mathcal{A}_{0}+\theta \mathcal{A}_{1}\right) u(t) d t=d W^{H}(t) \tag{3.1}
\end{equation*}
$$

with solution $u(t)=\sum_{j \geq 1} u_{j}(t) h_{j}$ given by (2.11); for simplicity, we assume that $u(0)=0$. Suppose that the processes $u_{1}(t), \ldots, u_{N}(t)$ can be observed for all $t \in$ $[0, T]$; as was mentioned in the Introduction, if $u$ is observable, then each $u_{k}$ can be computed. The problem is to estimate the parameter $\theta$ using these observations.

Recall the notation $\mu_{j}(\theta)=\rho_{j}+\nu_{j} \theta$, where $\rho_{j}$ and $\nu_{j}$ are the eigenvalues of $\mathcal{A}_{0}$ and $\mathcal{A}_{1}$, respectively. Then each $u_{j}$ is a fractional Ornstein-Uhlenbeck process satisfying

$$
\begin{equation*}
d u_{j}(t)=-\mu_{j}(\theta) u_{j}(t) d t+d w_{j}^{H}(t), \quad u_{j}(0)=0 \tag{3.2}
\end{equation*}
$$

and, because of the independence of $w_{j}^{H}$ for different $j$, the processes $u_{1}, \ldots, u_{N}$ are (statistically) independent.

Let $\Gamma$ denote the Gamma-function (see (2.12)). Following Kleptsyna and Le Breton [7], we introduce the notations

$$
\begin{gather*}
\kappa_{H}=2 H \Gamma\left(\frac{3}{2}-H\right) \Gamma\left(H+\frac{1}{2}\right), \quad k_{H}(t, s)=\kappa_{H}^{-1} s^{\frac{1}{2}-H}(t-s)^{\frac{1}{2}-H} ;  \tag{3.3}\\
\lambda_{H}=\frac{2 H \Gamma(3-2 H) \Gamma\left(H+\frac{1}{2}\right)}{\Gamma\left(\frac{3}{2}-H\right)}, \quad \mathbf{w}_{H}(t)=\lambda_{H}^{-1} t^{2-2 H} ; \tag{3.4}
\end{gather*}
$$

$$
\begin{gather*}
M_{j}^{H}(t)=\int_{0}^{t} k_{H}(t, s) d w_{j}^{H}(s), \quad Q_{j}(t)=\frac{d}{d \mathbf{w}_{H}(t)} \int_{0}^{t} k_{H}(t, s) u_{j}(s) d s  \tag{3.5}\\
Z_{j}(t)=\int_{0}^{t} k_{H}(t, s) d u_{j}(s) . \tag{3.6}
\end{gather*}
$$

In particular, $M^{H}$ is a martingale with the quadratic characteristic $\mathbf{w}_{H}$. By a Girsanov-type formula (see, for example, Kleptsyna et al. [8, Theorem 3]), there is a probability measure $\widetilde{\mathbb{P}}$ on $(\Omega, \mathcal{F})$ under which the distribution of the vector $\left(u_{1}, \ldots, u_{N}\right)$ is the same as the distribution of the vector $\left(w_{1}^{H}, \ldots, w_{N}^{H}\right)$ under the original measure $\mathbb{P}$, and

$$
\begin{equation*}
\frac{d \widetilde{\mathbb{P}}}{d \mathbb{P}}=\exp \left(\sum_{j=1}^{N} \mu_{j}(\theta) \int_{0}^{T} Q_{j}(s) d Z_{j}(s)+\sum_{j=1}^{N} \frac{\left|\mu_{j}(\theta)\right|^{2}}{2} \int_{0}^{T} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)\right) . \tag{3.7}
\end{equation*}
$$

Note that process $Z_{j}$ is a semi-martingale [7, Lemma 2.1], and so there is no stochastic integration with respect to fractional Brownian motion in (3.7): $\int_{0}^{T} \nu_{j} Q_{j}(s) d Z_{j}(s)$ is an Itô integral. Maximizing the density $d \widetilde{\mathbb{P}} / d \mathbb{P}$ with respect to $\theta$ gives the Maximum Likelihood Estimator (MLE) of $\theta$ :

$$
\begin{equation*}
\widehat{\theta}_{N}=-\frac{\sum_{j=1}^{N} \int_{0}^{T} \nu_{j} Q_{j}(s)\left(d Z_{j}(s)+\rho_{j} Q_{j}(s) d \mathbf{w}_{H}(s)\right)}{\sum_{j=1}^{N} \int_{0}^{T} \nu_{j}^{2} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)} . \tag{3.8}
\end{equation*}
$$

Under assumption that $u_{1}, \ldots, u_{N}$ are observable, the processes $Q_{j}, Z_{j}, j=$ $1, \ldots, N$ are also observable, and hence the estimator $\hat{\theta}_{N}$ defined above is observable. On the other hand, computing (3.8) requires the knowledge of $k_{H}$ and $\mathbf{w}_{H}$, and is therefore impossible without the knowledge of the Hurst parameter $H$.

Notice that, when $H=1 / 2$, we have $k_{H}=1, \mathbf{w}_{H}(s)=s, Q_{j}(s)=Z_{j}(s)=u_{j}(s)$, and (3.8) becomes

$$
\begin{equation*}
\widehat{\theta}_{N}=-\frac{\sum_{j=1}^{N} \int_{0}^{T} \nu_{j} u_{j}(s)\left(d u_{j}(s)+\rho_{j} u_{j}(s) d s\right)}{\sum_{j=1}^{N} \int_{0}^{T} \nu_{j}^{2} u_{j}^{2}(s) d u_{j}(s)}, \tag{3.9}
\end{equation*}
$$

which is the MLE from [3].
The following is the main result of the paper.
Theorem 3.1. Assume that Eq. (3.1) is diagonalizable and parabolic, and $\theta \in \Theta$. Then the following conditions are equivalent:
(1) $\sum_{j=J}^{\infty} \frac{\nu_{j}^{2}}{\mu_{j}(\theta)}=+\infty$;
(2) $\lim _{N \rightarrow \infty} \widehat{\theta}_{N}=\theta$ with probability one,
where $J=\min \left\{j: \mu_{i}(\theta)>0\right.$ for all $\left.i \geq j\right\}$.

Proof. Following Kleptsyna and Le Breton [7, Eq. (4.1)], we conclude that

$$
\begin{equation*}
\widehat{\theta}_{N}-\theta=-\frac{\sum_{j=1}^{N} \int_{0}^{T} \nu_{j} Q_{j}(s) d M_{j}^{H}(s)}{\sum_{j=1}^{N} \int_{0}^{T} \nu_{j}^{2} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)} \tag{3.12}
\end{equation*}
$$

Both the top and the bottom on the right-hand side of (3.12) are sums of independent random variables; moreover, it is known from [7, p. 242] that

$$
\begin{equation*}
\mathbb{E}\left(\int_{0}^{T} Q_{j}(s) d M_{j}^{H}(s)\right)^{2}=\mathbb{E} \int_{0}^{T} Q_{j}^{2}(s) d \mathbf{w}_{H}(s) d s \tag{3.13}
\end{equation*}
$$

From the expression for the Laplace transform of $\int_{0}^{T} Q_{j}^{2}(s) d \mathbf{w}_{H}(s) d s$ (see [7, Eq. (4.2)]) direct computations show that

$$
\begin{equation*}
\lim _{j \rightarrow \infty} \mu_{j}(\theta) \mathbb{E} \int_{0}^{T} Q_{j}^{2}(s) d \mathbf{w}_{H}(s) d s=\frac{T}{2}>0 \tag{3.14}
\end{equation*}
$$

and, with $\operatorname{Var}(\xi)$ denoting the variance of the random variable $\xi$,

$$
\begin{equation*}
\lim _{j \rightarrow \infty} \mu_{j}^{3}(\theta) \operatorname{Var}\left(\int_{0}^{T} Q_{j}^{2}(s) d \mathbf{w}_{H}(s) d s\right)=\frac{T}{2}>0 \tag{3.15}
\end{equation*}
$$

a detailed derivation of (3.14) and (3.15) is given in the Appendix, Lemmas A. 1 and A. 2 respectively.

We now see that if (3.10) does not hold, then, by (3.14), the series

$$
\sum_{j \geq 1} \int_{0}^{T} \nu_{j}^{2} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)
$$

converges with probability one, which, by (3.12), means that (3.11) cannot hold.
On the other hand, if (3.10) holds, then

$$
\begin{equation*}
\sum_{n \geq J} \frac{\nu_{n}^{2} \mu_{n}^{-1}}{\left(\sum_{j=1}^{n} \nu_{j}^{2} \mu_{j}^{-1}\right)^{2}}<\infty \tag{3.16}
\end{equation*}
$$

Indeed, setting $a_{n}=\nu_{n}^{2} \mu_{n}^{-1}$ and $A_{n}=\sum_{j=1}^{n} a_{j}$, we notice that

$$
\sum_{n \geq J} \frac{a_{n}}{A_{n}^{2}} \leq \sum_{n \geq J+1}\left(\frac{1}{A_{n-1}}-\frac{1}{A_{n}}\right)=\frac{1}{A_{J}}
$$

Then the strong law of large numbers, together with the observation

$$
\mathbb{E} \int_{0}^{T} Q_{j}(s) d M_{j}^{H}(s)=0, \quad j \geq 1
$$

implies

$$
\lim _{N \rightarrow \infty} \frac{\sum_{j=1}^{N} \int_{0}^{T} \nu_{j} Q_{j}(s) d M_{j}(s)}{\sum_{j=1}^{N} \mathbb{E} \int_{0}^{T} \nu_{j}^{2} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)}=0 \quad \text { with probability one. }
$$

Next, it follows from (3.16) and the parabolicity assumption that

$$
\begin{equation*}
\sum_{n \geq J} \frac{\nu_{n}^{4} \mu_{n}^{-3}}{\left(\sum_{j=J}^{n} \nu_{j}^{2} \mu_{j}^{-1}\right)^{2}}<\infty \tag{3.17}
\end{equation*}
$$

because $\left|\nu_{j} / \mu_{j}(\theta)\right|$ stays bounded for $j \geq J$. Then another application of the strong law of large numbers implies that

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \frac{\sum_{j=1}^{N} \int_{0}^{T} \nu_{j}^{2} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)}{\sum_{j=1}^{N} \mathbb{E} \int_{0}^{T} \nu_{j}^{2} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)}=1 \tag{3.18}
\end{equation*}
$$

with probability one, and (3.11) follows.
Corollary 3.1. If Eq. (3.1) is diagonalizable and parabolic and if (3.10) holds, then

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \sqrt{\sum_{j=J}^{N} \frac{\nu_{j}^{2}}{\mu_{j}(\theta)}}\left(\widehat{\theta}_{N}-\theta\right)=\zeta \tag{3.19}
\end{equation*}
$$

in distribution, where $\zeta$ is a Gaussian random variable with zero mean and variance $2 / T$.

Proof. Consider the sequence of random processes

$$
X_{N}(t)=\frac{\sum_{j=1}^{N} \int_{0}^{t} \nu_{j} Q_{j}(s) d M_{j}(s)}{\left(\sum_{j=1}^{N} \mathbb{E} \int_{0}^{T} \nu_{j}^{2} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)\right)^{1 / 2}}
$$

and let $X(t)=w(t) / T$, where $w$ is a standard Brownian motion. It is known from Norros et al. [12, Theorem 3.1] that each $M_{j}^{H}$ is a Gaussian martingale with independent increments and quadratic variation $\left\langle M_{j}^{H}\right\rangle(t)=\mathbf{w}_{H}(t)$. Therefore, each $X_{N}$ is a square-integrable martingale with quadratic variation

$$
\left\langle X_{N}\right\rangle(t)=\frac{\sum_{j=1}^{N} \int_{0}^{t} \nu_{j}^{2} Q_{j}^{2}(s) d \mathbf{w}_{H}(t)}{\sum_{j=1}^{N} \mathbb{E} \int_{0}^{T} \nu_{j}^{2} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)} .
$$

By (3.18),

$$
\lim _{N \rightarrow \infty}\left\langle X_{N}\right\rangle(T)=1=\langle X\rangle(T)
$$

By a limit theorem for martingales (see, for example, Jacod and Shiryaev [6, Theorem VIII.4.17]), $\lim _{N \rightarrow \infty} X_{N}(T)=X(T)$ in distribution. It remains to notice that, by (3.14),

$$
\lim _{N \rightarrow \infty} \frac{\sum_{j=1}^{N} \nu_{j}^{2} / \mu_{j}(\theta)}{\sum_{j=1}^{N} \mathbb{E} \int_{0}^{T} \nu_{j}^{2} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)}=\frac{2}{T}
$$

By Corollary 3.1, note that the rate of convergence of the estimators $\widehat{\theta}_{N}$ does not depend on the Hurst parameter $H$.

## 4. Examples

First, we discuss how the presence of the spatial covariance in the noise term affects the model.

Let us consider a more general equation

$$
d u(t)=\left(\mathcal{A}_{0}+\theta \mathcal{A}_{1}\right) u(t) d t+\mathcal{B} d W^{H}(t)
$$

where $\mathcal{B}$ is a linear operator, with $\mathcal{B B}^{\top}$ representing the spatial covariance of the noise process. If $\mathcal{B}^{-1}$ exists, the equation is reduced to (3.1) by considering $v=$ $\mathcal{B}^{-1} u, \tilde{\mathcal{A}}_{0}=\mathcal{B}^{-1} \mathcal{A}_{0} \mathcal{B}, \tilde{\mathcal{A}}_{1}=\mathcal{B}^{-1} \mathcal{A}_{1} \mathcal{B}:$

$$
d v(t)=\left(\tilde{\mathcal{A}}_{0}+\theta \tilde{\mathcal{A}}_{1}\right) v(t) d t+d W^{H}(t) ;
$$

to proceed, we certainly need to assume that this equation is diagonalizable and parabolic.

If $\mathcal{B}^{-1}$ does not exist, we have two possibilities:
(1) $\left(u_{0}, h_{i}\right)_{0}=0$ for every $i$ such that $\mathcal{B} h_{i}=0$. In this case, $u_{i}(t)=0$ for all $t>0$, so that we can factor out the kernel of $\mathcal{B}$ and reduce the problem to invertible $\mathcal{B}$.
(2) $\left(u_{0}, h_{i}\right)_{0} \neq 0$ for some $i$ such that $\mathcal{B} h_{i}=0$. In this case, $u_{i}(t)=u_{i}(0) e^{-\rho_{i} t-\nu_{i} \theta t}$ and $\theta$ is determined exactly from the observations of $u_{i}(t)$ :

$$
\theta=\frac{1}{\nu_{i}(t-s)} \ln \frac{u_{i}(s)}{u_{i}(t)}-\frac{\rho_{i}}{\nu_{i}}, \quad t \neq s
$$

Next, let us formulate condition (3.10) in terms of the orders of the operators in the equation. Let $\mathcal{A}_{0}, \mathcal{A}_{1}$ be differential or pseudo-differential operators, either on a smooth bounded domain in $\mathbb{R}^{d}$ or on a smooth compact $d$-dimensional manifold, and let $m_{0}, m_{1}$, be the orders of $\mathcal{A}_{0}, \mathcal{A}_{1}$ respectively, so that $2 m=\max \left(m_{0}, m_{1}\right)$. Then, under rather general conditions we have

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left|\nu_{j}\right| j^{m_{1} / d}=c_{1}, \quad \lim _{j \rightarrow \infty} \mu_{j}(\theta) j^{2 m / d}=c(\theta) \tag{4.1}
\end{equation*}
$$

for some positive numbers $c_{1}, c(\theta)$; see, for example, Il'in [5] or Safarov and Vassiliev [15].

If (4.1) holds, then condition (3.10) becomes

$$
\begin{equation*}
m_{1} \geq m-(d / 2) \tag{4.2}
\end{equation*}
$$

which, in the case $H=1 / 2$, was established by Huebner and Rozovskii [3]. On the other hand, Theorem 3.1 covers operators with more exotic eigenvalues, such as $\nu_{j}=j \ln j$ or $\nu_{j}=e^{j}$.

Note that, at least as long as $H \geq 1 / 2$, conditions (3.10) and (4.2) do not involve $H$.

Let us now look at some concrete examples.
(1) Consider the equation

$$
d u-\theta u_{x x} d t=d W^{H}, \quad 0<t<T, \quad x \in(0,1)
$$

with zero initial and boundary conditions. This equation was also studied by Prakasa Rao [14].

Clearly, $\nu_{j}=\pi^{2} j^{2}, \rho_{j}=0, \mu_{j}(\theta)=\theta \pi^{2} j^{2}$. We take $\Theta=\left(\theta_{0}, \theta_{1}\right)$ for some $0<\theta_{0}<\theta_{1}$. Then

$$
\begin{gathered}
u_{j}(t)=\int_{0}^{T} e^{-\theta \pi^{2} j^{2}(t-s)} d w_{j}^{H} \\
\widehat{\theta}_{N}=-\frac{\sum_{j=1}^{N} j^{2} \int_{0}^{T} Q_{j}(t) d Z_{j}(t)}{\sum_{j=1}^{N} \int_{0}^{T} \pi^{2} j^{4} Q_{j}^{2}(t) d \mathbf{w}_{H}(t)}, \\
\lim _{N \rightarrow \infty} N^{3 / 2}\left(\widehat{\theta}_{N}-\theta\right)=\mathcal{N}\left(0,6 \theta /\left(\pi^{2} T\right)\right),
\end{gathered}
$$

where $\mathcal{N}\left(a, \sigma^{2}\right)$ is a normal random variable with mean $a$ and variance $\sigma^{2}$, and the convergence is in distribution.
(2) Consider the equation

$$
d u-(\boldsymbol{\Delta} u+\theta u) d t=d W^{H}, \quad 0<t<T, x \in G \subset \mathbb{R}^{d}
$$

with zero initial and boundary conditions and $d \geq 2$. Denote by $\lambda_{j}, j \geq 1$ the eigenvalues of the Laplace operator $\boldsymbol{\Delta}$; recall that $\lambda_{j}<0$. Clearly, $\nu_{j}=-1, \rho_{j}=$ $-\lambda_{j}, \mu_{j}(\theta)=-\lambda_{j}-\theta$. We take $\Theta=\left(\theta_{0}, \theta_{1}\right)$ for some $\theta_{0}<\theta_{1}$; there are no restrictions on the sign of $\theta$. Then

$$
\begin{array}{r}
u_{j}(t)=\int_{0}^{T} e^{\left(\lambda_{j}+\theta\right)(t-s)} d w_{j}^{H}, \\
\widehat{\theta}_{N}=-\frac{\sum_{j=1}^{N} \int_{0}^{T} Q_{j}(t)\left(d Z_{j}(t)-\lambda_{j} Q_{j}(t) d \mathbf{w}_{H}(t)\right)}{\sum_{j=1}^{N} \int_{0}^{T} Q_{j}^{2}(t) d \mathbf{w}_{H}(t)}, \\
\lim _{N \rightarrow \infty} \Psi_{d}(N)\left(\widehat{\theta}_{N}-\theta\right)=\mathcal{N}\left(0, \sigma_{d}^{2}\right),
\end{array}
$$

where

$$
\Psi_{d}(N)=\left\{\begin{array}{ll}
\sqrt{\ln N}, & \text { if } d=2, \\
N^{(d-2) /(2 d)}, & \text { if } d>2,
\end{array} \quad \sigma_{d}^{2}= \begin{cases}\frac{2 c}{T}, & \text { if } d=2 \\
\frac{2 c}{T}\left(1-\frac{2}{d}\right), & \text { if } d>2\end{cases}\right.
$$

and $c$ is from (2.6).
If $d=1$, then (4.2) does not hold and $\widehat{\theta}_{N}$ is not a consistent estimator of $\theta$.
(3) Consider the equation

$$
d u+\left(\boldsymbol{\Delta}^{2} u+\theta \boldsymbol{\Delta} u\right) d t=d W^{H}, \quad 0<t<T, x \in G \subset \mathbb{R}^{d}
$$

with zero initial and boundary conditions:

$$
\left.u\right|_{t=0}=\left.u\right|_{\partial G}=\left.\boldsymbol{\Delta} u\right|_{\partial G}=0 .
$$

As before, denote by $\lambda_{j}$ the eigenvalues of the Laplacian $\boldsymbol{\Delta}$. Clearly, $\nu_{j}=-\lambda_{j}$, $\rho_{j}=\lambda_{j}^{2}, \mu_{j}(\theta)=\lambda_{j}^{2}+\theta \lambda_{j}$. We take $\Theta=\left(\theta_{0}, \theta_{1}\right)$ for some $\theta_{0}<\theta_{1}$; there are no restrictions on the sign of $\theta$. Then, for every $d \geq 1$,

$$
\begin{array}{r}
u_{j}(t)=\int_{0}^{T} e^{-\left(\lambda_{j}^{2}+\theta \lambda_{j}\right)(t-s)} d w_{j}^{H} \\
\widehat{\theta}_{N}=-\frac{\sum_{j=1}^{N} \int_{0}^{T} \lambda_{j} Q_{j}(t)\left(d Z_{j}(t)-\lambda_{j}^{2} Q_{j}(t) d \mathbf{w}_{H}(t)\right)}{\sum_{j=1}^{N} \int_{0}^{T} \lambda_{j}^{2} Q_{j}^{2}(t) d \mathbf{w}_{H}(t)} \\
\lim _{N \rightarrow \infty} \sqrt{N}\left(\widehat{\theta}_{N}-\theta\right)=\mathcal{N}(0,2 / T)
\end{array}
$$

## 5. Other Possible Estimators

The maximum likelihood estimator (3.8) has three features that are clearly attractive: consistency, asymptotic normality, and absence of stochastic integration with respect to fractional Brownian motion. On the other hand, actual implementation of (3.8) is problematic: when $H>1 / 2$, computing the processes $Q_{j}$ and $Z_{j}$ is certainly nontrivial. Estimator (3.9) is defined for all $H \geq 1 / 2$ and contains only the processes $u_{j}$, but, when $H>1 / 2$, is not an MLE and is even harder to implement because of the stochastic integral with respect to $u_{j}$.

With or without condition (3.10), a consistent estimator of $\theta$ is possible in the large time asymptotic: for every $j \geq 1$,

$$
\begin{equation*}
\lim _{T \rightarrow \infty} \frac{\int_{0}^{T} \nu_{j} Q_{j}(s)\left(d Z_{j}(s)+\rho_{j} Q_{j}(s) d \mathbf{w}_{H}(s)\right)}{\int_{0}^{T} \nu_{j}^{2} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)}=-\theta \tag{5.1}
\end{equation*}
$$

with probability one [7, Proposition 2.2]. For $H>1 / 2$, implementation of this estimator is essentially equivalent to the implementation of (3.8).

An alternative to (5.1) was suggested by Maslowski and Pospísil [11] using the ergodic properties of the OU process. Let us first illustrate the idea on a simple example.

If $a>0$ and $w=w(t)$ is a standard one-dimensional Brownian motion, then the OU process $d X=-a X(t) d t+d w(t)$ is ergodic and its unique invariant distribution is normal with zero mean and variance $(2 a)^{-1}$. In particular,

$$
\begin{equation*}
\lim _{T \rightarrow \infty} \frac{1}{T} \int_{0}^{T} X^{2}(t) d t=\frac{1}{2 a} \tag{5.2}
\end{equation*}
$$

with probability one, and so

$$
\begin{equation*}
\tilde{a}(T)=\frac{T}{2 \int_{0}^{T} X^{2}(t) d t} \tag{5.3}
\end{equation*}
$$

is a consistent estimator of $a$ in the long-time asymptotic. Note that the maximum likelihood estimator in this case is

$$
\begin{equation*}
\hat{a}(T)=-\frac{\int_{0}^{T} X(t) d X(t)}{\int_{0}^{T} X^{2}(s) d s} \tag{5.4}
\end{equation*}
$$

and is strongly consistent for every $a \in \mathbb{R}$ [10, Theorem 17.4].
Similarly, if $a>0$, then the fractional OU process

$$
\begin{equation*}
d X(t)=-a X(t) d t+d w^{H}(t), \quad X(0)=0 \tag{5.5}
\end{equation*}
$$

is Gaussian, and, by (2.12), converges in distribution, as $t \rightarrow \infty$, to the Gaussian random variable with zero mean and variance $c(H) a^{-2 H}$, where

$$
\begin{equation*}
c(H)=H(2 H-1) \Gamma(2 H-1) ; \tag{5.6}
\end{equation*}
$$

notice that, in the limit $H \searrow 1 / 2$, we recover the result for the usual OU process. Further investigation shows that, similar to (5.2),

$$
\lim _{T \rightarrow \infty} \frac{1}{T} \int_{0}^{T} X^{2}(s) d s=\frac{c(H)}{a^{2 H}}
$$

(see [11]). As a result, for every $j$ such that $\theta \nu_{j}+\rho_{j}>0$, we have

$$
\begin{equation*}
\lim _{T \rightarrow \infty} \frac{1}{T} \int_{0}^{T} u_{j}^{2}(t) d t=\frac{c(H)}{\left(\theta \nu_{j}+\rho_{j}\right)^{2 H}} \tag{5.7}
\end{equation*}
$$

with probability one. Under an additional assumption that $\nu_{j} \neq 0$, we get an estimator of $\theta$

$$
\begin{equation*}
\tilde{\theta}^{(j)}(T)=\frac{1}{\nu_{j}}\left(\frac{c(H) T}{\int_{0}^{T} u_{j}^{2}(t) d t}\right)^{\frac{1}{2 H}}-\frac{\rho_{j}}{\nu_{j}} . \tag{5.8}
\end{equation*}
$$

This estimator is strongly consistent in the long time asymptotic: $\lim _{T \rightarrow \infty} \mid \tilde{\theta}^{(j)}(T)-$ $\theta \mid=0$ with probability one [11, Theorem 5.2]. While not a maximum likelihood estimator, (5.8) is easier to implement computationally than (3.8). If, in Theorem 2.1, we have $\mathcal{A}_{0}=0$ and $\nu_{j}>0$, then a version of (5.9) exists using all the Fourier coefficients $u_{j}, j \geq 1$ :

$$
\begin{equation*}
\tilde{\theta}(T)=\left(\frac{c(H) T \sum_{j=1}^{\infty} \nu_{j}^{-2 H}}{\sum_{j=1}^{\infty} \int_{0}^{T} u_{j}^{2}(t) d t}\right)^{\frac{1}{2 H}} \tag{5.9}
\end{equation*}
$$

see [11, Theorem 5.2].
An interesting open question related to both (3.8) and (5.8), (5.9) is how to combine estimation of $\theta$ with estimation of $H$.

## Appendix A

Below, we prove equalities (3.14) and (3.15).
Lemma A.1. For every $\theta \in \Theta$ and $H \in[1 / 2,1)$,

$$
\lim _{j \rightarrow \infty} \mu_{j}(\theta) \mathbb{E} \int_{0}^{T} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)=\frac{T}{2}
$$

Proof. Denote by $\Psi_{T}^{H}\left(a, \mu_{j}\right)$ the Laplace transform of $\int_{0}^{T} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)$, namely

$$
\begin{equation*}
\Psi_{T}^{H}\left(a, \mu_{j}(\theta)\right)=\mathbb{E} \exp \left\{-a \int_{0}^{T} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)\right\}, \quad a>0 \tag{A.1}
\end{equation*}
$$

We will use the expression for $\Psi_{T}^{H}$ from [7, p. 242], and write it as follows

$$
\Psi_{T}^{H}\left(a, \mu_{j}\right)=\alpha e^{\frac{\left(\mu_{j}-\alpha\right) T}{2}}\left[\Delta_{T}^{H}\left(\mu_{j}, \alpha\right)\right]^{-\frac{1}{2}}
$$

where $\mu_{j}=\mu_{j}(\theta), \alpha:=\sqrt{\mu_{j}^{2}+2 a}$,

$$
\begin{aligned}
\Delta_{T}^{H}\left(\mu_{j}, \alpha\right)= & \frac{\pi \alpha T e^{-\alpha T}\left(\alpha^{2}-\mu_{j}^{2}\right)}{4 \sin (\pi H)} I_{-H}\left(\frac{\alpha T}{2}\right) I_{H-1}\left(\frac{\alpha T}{2}\right) \\
& +e^{-\alpha T}\left[\alpha \sinh \left(\frac{\alpha T}{2}\right)+\mu_{j} \cosh \left(\frac{\alpha T}{2}\right)\right]^{2}
\end{aligned}
$$

and $I_{p}$ is the modified Bessel function of the first kind and order $p$.
Note that

$$
\mathbb{E} \int_{0}^{T} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)=-\left.\frac{\partial \Psi_{T}^{H}\left(a, \mu_{j}\right)}{\partial a}\right|_{a=0}
$$

Direct evaluations (for example, using Mathematica computer algebra system) give

$$
\left.\frac{\partial \Psi_{T}^{H}\left(a, \mu_{j}\right)}{\partial a}\right|_{a=0}=\frac{2+2 e^{\mu_{j} T}\left(1-\mu_{j} T\right)-\mu_{j} \pi T I_{H-1}\left(\frac{\mu_{j} T}{2}\right) I_{-H}\left(\frac{\mu_{j} T}{2}\right) \csc (H \pi)}{4 \mu_{j}^{2} e^{\mu_{j} T}}
$$

where $\csc (x)=1 / \sin (x)$. By combining formulas (6.106), (6.155) and (6.162) in [1], we conclude that, for all $p \in(-1,1), p \neq 0$, we have $I_{p}(x) \sim e^{x} / \sqrt{2 \pi x}, x \rightarrow \infty$, i.e.

$$
\begin{equation*}
\lim _{x \rightarrow+\infty} \sqrt{2 \pi x} e^{-x} I_{p}(x)=1 \tag{A.2}
\end{equation*}
$$

Therefore

$$
\begin{aligned}
\left.\frac{\partial \Psi_{T}^{H}\left(a, \mu_{j}\right)}{\partial a}\right|_{a=0} \sim & \frac{2+2 e^{\mu_{j} T}\left(1-\mu_{j} T\right)-e^{\mu_{j} T} \csc (H \pi)}{4 \mu_{j}^{2} e^{\mu_{j} T}} \sim-\frac{T}{2 \mu_{j}}, \quad j \rightarrow \infty \\
& \left.\lim _{j \rightarrow \infty} \mu_{j} \frac{\partial \Psi_{T}^{H}\left(a, \mu_{j}\right)}{\partial a}\right|_{a=0}=-\frac{T}{2}
\end{aligned}
$$

and the lemma is proved.

Lemma A.2. For every $\theta \in \Theta$ and $H \in[1 / 2,1)$

$$
\lim _{j \rightarrow \infty} \mu_{j}^{3}(\theta) \operatorname{Var}\left(\int_{0}^{T} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)\right)=\frac{T}{2} .
$$

Proof. Note that

$$
\begin{equation*}
\mathbf{V}:=\operatorname{Var}\left(\int_{0}^{T} Q_{j}^{2}(s) d \mathbf{w}_{H}(s)\right)=\left[\frac{\partial^{2} \Psi_{T}^{H}\left(a, \mu_{j}\right)}{\partial a^{2}}-\left(\frac{\partial \Psi_{T}^{H}\left(a, \mu_{j}\right)}{\partial a}\right)^{2}\right]_{a=0} \tag{A.3}
\end{equation*}
$$

with $\Psi_{T}^{H}$ from (A.1). Direct evaluation of the right-hand side of (A.3) (for example, using Mathematica computer algebra system) gives

$$
\begin{aligned}
\mathbf{V}= & \frac{1}{8 \mu_{j}^{4} e^{2 T \mu_{j}}}\left(2-8 e^{\mu_{j} T}\left(1+\mu_{j} T\right)+2 e^{2 \mu_{j} T}\left(-5+2 \mu_{j} T\right)\right. \\
& +\pi \mu_{j} T \csc (\pi H)\left[-2 e^{\mu_{j} T} \mu_{j} T I_{1-H}\left(\frac{\mu_{j} T}{2}\right) I_{H-1}\left(\frac{\mu_{j} T}{2}\right)\right. \\
& +I_{-H}\left(\frac{\mu_{j} T}{2}\right)\left\{4\left(-1+e^{\mu_{j} T}\left(1+\mu_{j} T\right)\right) I_{H-1}\left(\frac{\mu_{j} T}{2}\right)\right. \\
& \left.\left.\left.-2 e^{\mu_{j} T} \mu_{j} T I_{H}\left(\frac{\mu_{j} T}{2}\right)+\pi \mu_{j} T I_{H-1}^{2}\left(\frac{\mu_{j} T}{2}\right) I_{-H}\left(\frac{\mu_{j} T}{2}\right) \csc (H \pi)\right\}\right]\right)
\end{aligned}
$$

where $\csc (x)=1 / \sin (x)$ and $I_{p}$ is the modified Bessel function of the first kind and order $p$.

Using (A.2), we conclude that

$$
\begin{aligned}
\lim _{j \rightarrow \infty} \mu_{j}^{3}(\theta) \mathbf{V}= & \lim _{j \rightarrow \infty} \mu_{j}^{3}\left(\frac{-10+4 \csc (H \pi)+\csc ^{2}(H \pi)}{8 \mu_{j}^{4}}+\frac{1}{4 \mu_{j}^{4} e^{2 \mu_{j} T}}\right. \\
& \left.-\frac{\csc (H \pi)+2+2 \mu_{j} T}{2 \mu_{j}^{4} e^{\mu_{j} T}}+\frac{T}{2 \mu_{j}^{3}}\right)=\frac{T}{2}
\end{aligned}
$$

and complete the proof of the lemma.

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