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# KOLMOGOROV EQUATION AND LARGE-TIME BEHAVIOUR FOR FRACTIONAL BROWNIAN MOTION DRIVEN LINEAR SDE'S* 

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Abstract. We consider a stochastic process $X_{t}^{x}$ which solves an equation

$$
\mathrm{d} X_{t}^{x}=A X_{t}^{x} \mathrm{~d} t+\Phi \mathrm{d} B_{t}^{H}, \quad X_{0}^{x}=x
$$

where $A$ and $\Phi$ are real matrices and $B^{H}$ is a fractional Brownian motion with Hurst parameter $H \in(1 / 2,1)$. The Kolmogorov backward equation for the function $u(t, x)=$ $\mathbb{E} f\left(X_{t}^{x}\right)$ is derived and exponential convergence of probability distributions of solutions to the limit measure is established.

Keywords: fractional Brownian motion, Kolmogorov backwards equation, linear stochastic equation

MSC 2000: 60H05, 60H10

## 1. Introduction

Fractional Brownian motion (fBm) is a family of Gaussian processes indexed by a parameter $H \in(0,1)$ called the Hurst parameter. The first result on this type of processes goes back to Kolmogorov [13], basic properties of (fBm) have been proved in the pioneering work of Mandelbrot and Van Ness [15]. Hurst [10], [11] used these processes to describe the long term storage capacity of reservoirs along the Nile river, which demonstrated their usefulness as a model for physical phenomena. Later these processes have been used to describe some economic data and most recently to model telecommunication traffic.

[^0]In recent years numerous papers on (fBm)-based stochastic calculus have appeared (e.g. [4] and [1], [3], [5], [6], [7], [20], [21]). Difficulties are coming from the fact that $(\mathrm{fBm})$ is not a semimartingal, so the classical concept of stochastic integral may not be used. On the other hand, the paths of ( fBm ) do not have bounded variation, hence the classical concept of Lebesgue-Stieltjes integral is useless as well and it is necessary to develop a stochastic calculus. Very recently stochastic differential equations driven by (fBm) have been studied, see e.g. [2], [12], [14], [17] in a finite dimensional state space and [7], [8], [16] in infinite dimensions.

In the present paper we are concerned with a linear stochastic differential equation driven by (fBm). We derive the backward Kolmogorov equation (BKE) associated with the linear stochastic equation. Although our stochastic equation is autonomous, the coefficients of the (BKE) are time-dependent, which reflects the non-Markovian character of the equation. Note that the PDE for the transition density of (fBm) has been derived in [19].

The definition and basic properties of a stochastic integral of a deterministic function with respect to the ( fBm ) and the formula for integration by parts are in Section 2 of the paper. The main results of the paper are Propositions 3.3, 3.4 and 4.2 where the parabolic equation for $\mathbb{E} f\left(X_{t}^{x}\right)$ is derived and conditions for exponential convergence in total variation of probability distributions of $X_{t}^{x}$ to the limit distribution are established. Note that in [7] strong convergence (but not the speed of convergence) has been shown in an infinite dimensional state space.

## 2. Preliminaries

Let $(\Omega, \mathcal{F}, \mathscr{P})$ be a probability space.
Definition 2.1. A fractional Brownian motion $\left(B^{H}(t), t \geqslant 0\right)$ with Hurst parameter $H \in(0,1)$ on $(\Omega, \mathcal{F}, \mathscr{P})$ is an $\mathbb{R}^{n}$-valued Gaussian process which satisfies
(i) $\mathbb{E} B^{H}(t)=0$ for all $t \in \mathbb{R}_{+}$,
(ii) $\mathbb{E} B^{H}(s) B^{H}(t)^{\prime}=1 / 2\left(s^{2 H}+t^{2 H}-|t-s|^{2 H}\right) I$, for all $s, t \in \mathbb{R}_{+}$, where $I$ is the identity matrix,
(iii) $\left(B^{H}(t), t \geqslant 0\right)$ has continuous sample paths $\mathscr{P}$-a.s.

Since the function $F^{H}(s, t)=1 / 2\left(s^{2 H}+t^{2 H}-|t-s|^{2 H}\right.$ is positive definite for $H \in(0,1)$ there exists a Gaussian process with this covariance function due to the Daniell-Kolmogorov theorem. Continuous modification of this process follows from the Kolmogorov-Tschentsov theorem. From the condition (ii) it is clear that the components of this process are independent. Conditions (i) and (ii) easily imply $B^{H}(0)=0 \mathscr{P}_{\text {-a.s. The }}(\mathrm{fBm}) B^{H}$ with $H=1 / 2$ is the usual Wiener process. In the rest of the paper we assume that $H>1 / 2$.

Now, following [7], we define a stochastic integral with respect to the (fBm), with deterministic integrand. For this purpose we use

Lemma 2.2. If $p>1 / H$, then for $\varphi \in L^{p}(0, T ; \mathbb{R})$ the inequality

$$
\int_{0}^{T} \int_{0}^{T} \varphi(u) \varphi(v) \psi(u-v) \mathrm{d} u \mathrm{~d} v \leqslant C|\varphi|_{L^{p}(0, T ; \mathbb{R})}^{2}
$$

holds for some $C>0$ which depends only on $T$ and $p$, where

$$
\begin{equation*}
\psi(u)=H(2 H-1)|u|^{2 H-2} . \tag{1}
\end{equation*}
$$

This lemma is an easy consequence of the Young inequality ([18], Proposition 12.25).

First we define the stochastic integral

$$
\begin{equation*}
\int_{0}^{T} f \mathrm{~d} \beta^{H} \tag{2}
\end{equation*}
$$

for $f \in L^{p}\left(0, T ; \mathbb{R}^{n}\right)$ and a scalar $(\mathrm{fBm})\left(\beta^{H}(t), t \in[0, T]\right)$. Let $\mathcal{E}$ be the family of all $\mathbb{R}^{n}$-valued step functions, i.e.

$$
\begin{aligned}
\mathcal{E}=\left\{f: f(s)=\sum_{i=1}^{n-1} f_{i} 1_{\left[t_{i}, t_{i+1}\right)}(s),\right. & 0=t_{1}<t_{2}<\ldots<t_{n}=T \quad \text { and } \\
& \left.f_{i} \in \mathbb{R}^{n} \quad \text { for } i \in\{1, \ldots, n-1\}\right\} .
\end{aligned}
$$

For $f \in \mathcal{E}$ we define the stochastic integral (2) as

$$
\int_{0}^{T} f \mathrm{~d} \beta^{H}:=\sum_{i=1}^{n-1} f_{i}\left(\beta^{H}\left(t_{i+1}\right)-\beta^{H}\left(t_{i}\right)\right)
$$

The mean of this Gaussian random variable is zero and the covariance matrix is

$$
\begin{align*}
\mathbb{E}\left(\int_{0}^{T} f \mathrm{~d} \beta^{H}\right. & \left.\int_{0}^{T} f^{\prime} \mathrm{d} \beta^{H}\right)  \tag{3}\\
= & \sum_{i=0}^{n-1} \sum_{j=0}^{n-1} f_{i} f_{j}^{\prime} \mathbb{E}\left(\beta^{H}\left(t_{i+1}\right)-\beta^{H}\left(t_{i}\right)\right)\left(\beta^{H}\left(t_{j+1}\right)-\beta^{H}\left(t_{j}\right)\right) \\
= & \sum_{i=0}^{n-1} \sum_{j=0}^{n-1} f_{i} f_{j}^{\prime} \mathbb{E}\left(\beta^{H}\left(t_{i+1}\right) \beta^{H}\left(t_{j+1}\right)-\beta^{H}\left(t_{i+1}\right) \beta^{H}\left(t_{j}\right)\right. \\
& \left.\quad-\beta^{H}\left(t_{i}\right) \beta^{H}\left(t_{j+1}\right)+\beta^{H}\left(t_{i}\right) \beta\left(t_{j}\right)\right) .
\end{align*}
$$

It is easy to verify that the last term is equal to

$$
\begin{equation*}
\sum_{i=0}^{n-1} \sum_{j=0}^{n-1} \int_{t_{i}}^{t_{i+1}} \int_{t_{j}}^{t_{j+1}} f_{i} f_{j}^{\prime} \psi(u-v) \mathrm{d} u \mathrm{~d} v=\int_{0}^{T} \int_{0}^{T} f(u) f(v)^{\prime} \psi(u-v) \mathrm{d} u \mathrm{~d} v \tag{4}
\end{equation*}
$$

By Lemma 2.2 it follows that

$$
\begin{align*}
\left\|\int_{0}^{T} f \mathrm{~d} \beta^{H}\right\|_{L^{2}\left(\Omega ; \mathbb{R}^{n}\right)}^{2} & =\int_{0}^{T} \int_{0}^{T}\langle f(u), f(v)\rangle \psi(u-v) \mathrm{d} u \mathrm{~d} v  \tag{5}\\
& \leqslant C\left(\int_{0}^{T}|f(s)|^{p} \mathrm{~d} s\right)^{2 / p}=C|f|_{L^{p}\left(0, T ; \mathbb{R}^{n}\right)}^{2}
\end{align*}
$$

Because $\mathcal{E}$ is dense in $L^{p}\left(0, T ; \mathbb{R}^{n}\right)$ there is a sequence $\left\{f_{k}\right\} \subset \mathcal{E}$ such that $\mid f_{k}-$ $\left.f\right|_{L^{p}\left(0, T ; \mathbb{R}^{n}\right)} \rightarrow 0$ for $f \in L^{p}\left(0, T ; \mathbb{R}^{n}\right)$. This is a Cauchy sequence in $L^{p}\left(0, T ; \mathbb{R}^{n}\right)$. It follows from inequality (5) that also the sequence of the corresponding stochastic integrals is a Cauchy sequence in the complete space $L^{2}\left(\Omega ; \mathbb{R}^{n}\right)$ and therefore $\int_{0}^{T} f_{k} \mathrm{~d} \beta^{H}$ converge to some centered, Gaussian random vector which we denote by $\int_{0}^{T} f \mathrm{~d} \beta^{H}$.

Now we will define the stochastic integral with respect to an $\mathbb{R}^{n}$-valued (fBm)

$$
\begin{equation*}
\int_{0}^{T} G \mathrm{~d} B^{H} \tag{6}
\end{equation*}
$$

where $G:[0, T] \rightarrow \mathbb{R}^{n \times n}$ is a non-random real matrix. We assume that the matrix $G$ satisfies

$$
\begin{equation*}
G(\cdot) \in L^{p}\left(0, T ; \mathbb{R}^{n \times n}\right) \tag{7}
\end{equation*}
$$

for some $p>1 / H$. We define the stochastic integral (6) as

$$
I(G ; T):=\int_{0}^{T} G \mathrm{~d} B^{H}:=\sum_{i=1}^{n} \int_{0}^{T} G_{i} \mathrm{~d} \beta_{i}^{H}
$$

where $\beta_{i}^{H}(\cdot), i=1, \ldots, n$, are components of the $(\mathrm{fBm}) B^{H}(\cdot)$ and $G_{i}$ are the columns of $G$. Independence of random variables $\int_{0}^{T} G_{i} \mathrm{~d} \beta_{i}^{H}$, conditions (7) and (5) imply

$$
\begin{aligned}
\mathbb{E}|I(G, T)|^{2} & =\sum_{i=1}^{n}\left|\int_{0}^{T} G_{i} \mathrm{~d} \beta^{H}\right|^{2} \\
& =\sum_{i=1}^{n} \int_{0}^{T} \int_{0}^{T}\left\langle G(s)_{i}, G(r)_{i}\right\rangle \psi(s-r) \mathrm{d} s \mathrm{~d} r<\infty
\end{aligned}
$$

therefore the integral (6) is an $\mathbb{R}^{n}$-valued Gaussian variable. Its probability law is given in the next proposition.

Proposition 2.3. If $G:[0, T] \rightarrow \mathbb{R}^{n \times n}$ satisfies condition (7), then $I(G, T)$ is an $\mathbb{R}^{n}$-valued Gaussian random variable with zero mean and the covariance matrix

$$
\begin{equation*}
\mathbb{E}\left[I(G, T) I(G, T)^{\prime}\right]=\int_{0}^{T} \int_{0}^{T} G(u) G(v)^{\prime} \psi(u-v) \mathrm{d} u \mathrm{~d} v \tag{8}
\end{equation*}
$$

where $\psi$ is given by (1).
Proof. $I(G, T)$ is obviously a Gaussian random variable with zero mean. Let $C$ be its covariance matrix. We have

$$
\begin{aligned}
\langle C x, y\rangle & =\mathbb{E}[\langle x, I(G, T)\rangle\langle y, I(G, T)\rangle] \\
& =\mathbb{E}\left[\left\langle x, \sum_{i=1}^{n} \int_{0}^{T} G_{i} \mathrm{~d} \beta_{i}^{H}\right\rangle\left\langle y, \sum_{j=1}^{n} \int_{0}^{T} G_{j} \mathrm{~d} \beta_{j}^{H}\right\rangle\right] \\
& =\mathbb{E}\left[\sum_{i=1}^{n}\left\langle x, \int_{0}^{T} G_{i} \mathrm{~d} \beta_{i}^{H}\right\rangle\left\langle y, \int_{0}^{T} G_{i} \mathrm{~d} \beta_{i}^{H}\right\rangle\right]
\end{aligned}
$$

for all $x, y \in \mathbb{R}^{n}$. For arbitrary $f \in L^{p}\left(0, T ; \mathbb{R}^{n}\right)$ we have

$$
\left\langle x, \int_{0}^{T} f \mathrm{~d} \beta_{i}^{H}\right\rangle:=\int_{0}^{T}\langle x, f\rangle \mathrm{d} \beta_{i}^{H} \quad \text { for } i=1, \ldots, n \text { and } x \in \mathbb{R}^{n}
$$

This is true for step functions and therefore for every function from $L^{p}\left(0, T ; \mathbb{R}^{n}\right)$. Consequently,

$$
\begin{aligned}
\langle C x, y\rangle & =\sum_{i=1}^{n} \mathbb{E}\left[\int_{0}^{T}\left\langle x, G_{i}\right\rangle \mathrm{d} \beta_{i}^{H} \int_{0}^{T}\left\langle y, G_{i}\right\rangle \mathrm{d} \beta_{i}^{H}\right] \\
& =\sum_{i=1}^{n} \int_{0}^{T} \int_{0}^{T}\left\langle x, G(r)_{i}\right\rangle\left\langle y, G(s)_{i}\right\rangle \psi(r-s) \mathrm{d} r \mathrm{~d} s \\
& =\int_{0}^{T} \int_{0}^{T}\left\langle G(s) G(r)^{\prime} x, y\right\rangle \psi(r-s) \mathrm{d} r \mathrm{~d} s
\end{aligned}
$$

and hence $C=\int_{0}^{T} \int_{0}^{T} G(u) G(v)^{\prime} \psi(u-v) \mathrm{d} u \mathrm{~d} v$.
Now it will be shown that the integration by parts is valid. Let $B V(0, T)$ denote the space of functions with bounded variation on $[0, T]$.

Proposition 2.4. Let $\varphi \in B V(0, T)$ and assume that $\beta^{H}(t)$ is an $\mathbb{R}$-valued ( $f B m$ ). Then

$$
\begin{equation*}
\int_{0}^{T} \varphi(t) \mathrm{d} \beta^{H}(t)=\varphi(T) \beta^{H}(T)-\int_{0}^{T} \beta^{H}(t) \mathrm{d} \varphi(t) \quad \mathscr{P} \text {-a.s. } \tag{9}
\end{equation*}
$$

where the second term on the right-hand side is the Stieltjes integral.
Proof. Since $\varphi \in B V(0, T)$ there exist nondecreasing functions $\varphi^{+}, \varphi^{-}$: $[0, T] \rightarrow \mathbb{R}$ such that $\varphi=\varphi^{+}-\varphi^{-}$. We will prove the equality (9) for $\varphi^{+}$(the proof for $\varphi^{-}$is analogous). Let $\Delta_{n}(T):=\left\{0=t_{0}^{n}<t_{1}^{n}<\ldots<t_{k_{n}}^{n}=T, k_{n} \in \mathbb{N}\right\}$, $n \in \mathbb{N}$, be a partition of the interval $[0, T]$ such that

$$
\begin{equation*}
\left|\Delta_{n}(T)\right|:=\max _{0 \leqslant k \leqslant k_{n}}\left|t_{k+1}^{n}-t_{k}^{n}\right| \rightarrow 0 \quad \text { for } n \rightarrow \infty \tag{10}
\end{equation*}
$$

If we define a sequence of the step functions

$$
\varphi_{n}^{+}(t)=\sum_{k=0}^{k_{n}-1} \varphi^{+}\left(t_{k}^{n}\right) 1_{\left[t_{k}^{n}, t_{k+1}^{n}\right)}(t) \quad \text { for } t \in[0, T), \quad \varphi_{n}^{+}(T)=\varphi(T) \text { for } n \in \mathbb{N},
$$

we obtain

$$
\begin{aligned}
& \sum_{k=0}^{k_{n}} \varphi^{+}\left(t_{k}^{n}\right)\left(\beta^{H}\left(t_{k+1}^{n}\right)-\beta^{H}\left(t_{k}^{n}\right)\right) \\
& \quad=\sum_{k=0}^{k_{n}}\left(\varphi^{+}\left(t_{k+1}^{n}\right) \beta^{H}\left(t_{k+1}^{n}\right)-\varphi^{+}\left(t_{k}^{n}\right) \beta^{H}\left(t_{k}^{n}\right)-\beta^{H}\left(t_{k+1}^{n}\right)\left(\varphi^{+}\left(t_{k+1}^{n}\right)-\varphi^{+}\left(t_{k}^{n}\right)\right)\right)
\end{aligned}
$$

and hence

$$
\begin{equation*}
\int_{0}^{T} \varphi_{n}^{+}(t) \mathrm{d} \beta^{H}(t)=\varphi^{+}(T) \beta^{H}(T)-\sum_{k=0}^{k_{n}-1} \beta^{H}\left(t_{k+1}^{n}\right)\left(\varphi^{+}\left(t_{k+1}^{n}\right)-\varphi^{+}\left(t_{k}^{n}\right)\right) \tag{11}
\end{equation*}
$$

For every $\omega \in \Omega$ the function $\beta^{H}(\omega, \cdot)$ is continuous, therefore for $n \rightarrow \infty$ the limit on the right-hand side exists and is equal to the Stieltjes integral

$$
\int_{0}^{T} \beta^{H}(t) \mathrm{d} \varphi^{+}(t)
$$

Functions $\varphi_{n}^{+}$converge to $\varphi^{+}$in $L^{p}(0, T ; \mathbb{R})$, because $1<1 / H<p$ and

$$
\begin{aligned}
\int_{0}^{T}\left|\varphi^{+}(t)-\varphi_{n}^{+}(t)\right|^{p} \mathrm{~d} t & =\sum_{k=0}^{k_{n}-1} \int_{t_{k}^{n}}^{t_{k+1}^{n}}\left|\varphi^{+}(t)-\varphi_{n}^{+}(t)\right|^{p} \mathrm{~d} t \\
& \leqslant \sum_{k=0}^{k_{n}-1}\left|t_{k+1}^{n}-t_{k}^{n}\right|\left(\varphi^{+}\left(t_{k+1}^{n}\right)-\varphi^{+}\left(t_{k}^{n}\right)\right)^{p} \\
& \leqslant\left|\Delta_{n}(T)\right|\left(\varphi^{+}(T)-\varphi^{+}(0)\right)^{p}
\end{aligned}
$$

Consequently, we can find a subsequence $\Delta_{n_{m}}(T)$ such that

$$
\int_{0}^{T} \varphi_{n_{m}}^{+}(t) \mathrm{d} \beta^{H}(t) \rightarrow \int_{0}^{T} \varphi^{+}(t) \mathrm{d} \beta^{H}(t) \quad \text { for } m \rightarrow \infty \quad \mathscr{P} \text {-a.s. }
$$

From the linearity of the stochastic integral we can conclude that $\varphi \in L^{p}(0, T ; \mathbb{R})$ and

$$
\begin{aligned}
\int_{0}^{T} \varphi(t) \mathrm{d} \beta^{H}(t) & :=\int_{0}^{T} \varphi^{+}(t) \mathrm{d} \beta^{H}(t)-\int_{0}^{T} \varphi^{-}(t) \mathrm{d} \beta^{H}(t) \\
& :=\varphi(T) \beta^{H}(T)-\left(\int_{0}^{T} \beta^{H}(t) \mathrm{d} \varphi^{+}(t)-\int_{0}^{T} \beta^{H}(t) \mathrm{d} \varphi^{-}(t)\right) \\
& =\varphi(T) \beta^{H}(T)-\int_{0}^{T} \beta^{H}(t) \mathrm{d} \varphi(t)
\end{aligned}
$$

## 3. The linear equation

In this section the linear stochastic differential equation

$$
\begin{equation*}
\mathrm{d} X_{t}=A X_{t} \mathrm{~d} t+\Phi \mathrm{d} B_{t}^{H}, \quad X_{0}=x_{0} \tag{12}
\end{equation*}
$$

is studied where $A \in \mathbb{R}^{n \times n}, \Phi \in \mathbb{R}^{n \times m}, x_{0} \in \mathbb{R}^{n}$ and $B^{H}$ is an $\mathbb{R}^{m}$-valued ( fBm ) defined on the probability space $(\Omega, \mathcal{F}, \mathscr{P})$.
$A$ strong solution of the equation (12) is defined as a continuous process $\left(X_{t}, t \geqslant 0\right)$ satisfying

$$
\begin{equation*}
X_{t}=x_{0}+\int_{0}^{t} A X_{s} \mathrm{~d} s+\Phi B_{t}^{H}, \quad t \geqslant 0 \quad \mathscr{P} \text {-a.s.. } \tag{13}
\end{equation*}
$$

The strong solution to the equation (12) exists and is given by the so called "mild formula":

$$
\begin{equation*}
X_{t}=\mathrm{e}^{A t} x_{0}+\int_{0}^{t} \mathrm{e}^{A(t-s)} \Phi \mathrm{d} B_{s}^{H}, \quad t \geqslant 0 \tag{14}
\end{equation*}
$$

As is well known, formula (14) for $H=1 / 2$ defines the Ornstein-Uhlenbeck process which solves our equation (12) with $H=1 / 2$ ( $B^{1 / 2}$ is the Wiener process).

The process defined by (14) is Gaussian with the mean $\mathrm{e}^{A t} x_{0}$ and the covariance matrix

$$
\begin{equation*}
Q_{t}=\int_{0}^{t} \int_{0}^{t} \mathrm{e}^{A u} \Phi \Phi^{\prime} \mathrm{e}^{A^{\prime} v} \psi(u-v) \mathrm{d} u \mathrm{~d} v \tag{15}
\end{equation*}
$$

for $t \geqslant 0$ as follows from Proposition 2.3.

Proposition 3.1. The strong solution to the equation (12) is unique $\mathscr{P}_{\text {-a.s. and }}$ is given by the formula (14).

Proof. First we will prove that the process (14) has a continuous modification. For $n \in \mathbb{N}$ and $t, s \in[n, n+1], t>s$, we have

$$
\begin{aligned}
\mathbb{E}\left|X_{t}-X_{s}\right|^{2} & \leqslant 2\left|\mathrm{e}^{A t} x_{0}-\mathrm{e}^{A s} x_{0}\right|^{2}+2 \mathbb{E}\left|\int_{0}^{t} \mathrm{e}^{A(t-u)} \Phi \mathrm{d} B_{u}^{H}-\int_{0}^{s} \mathrm{e}^{A(s-u)} \Phi \mathrm{d} B_{u}^{H}\right|^{2} \\
& =J_{1}+J_{2}
\end{aligned}
$$

There exist constants $K_{1}$ and $K_{2}$ such that

$$
\begin{aligned}
J_{1} \leqslant & 2\left(\left|x_{0}\right||A| \mathrm{e}^{|A| T}(t-s)\right)^{2} \leqslant K_{1}(t-s)^{2 H}, \\
J_{2}= & 2 \mathbb{E}\left|\int_{0}^{s}\left(\mathrm{e}^{A(t-u)}-\mathrm{e}^{A(s-u)} \Phi\right) \mathrm{d} B_{u}^{H}+\int_{s}^{t} \mathrm{e}^{A(t-s)} \Phi \mathrm{d} B_{u}^{H}\right|^{2} \\
\leqslant & 4 \mathbb{E}\left|\left(\mathrm{e}^{A(t-s)}-\mathrm{I}\right) \int_{0}^{s} \mathrm{e}^{A(s-u)} \Phi \mathrm{d} B_{u}^{H}\right|^{2}+4 \mathbb{E}\left|\int_{s}^{t} \mathrm{e}^{A(t-u)} \Phi \mathrm{d} B_{u}^{H}\right|^{2} \\
\leqslant & 4\left(|A| \mathrm{e}^{|A|}(t-s)\right)^{2} \operatorname{Tr} \int_{0}^{s} \int_{0}^{s} \mathrm{e}^{A u} \Phi \Phi^{\prime} \mathrm{e}^{A^{\prime} v} \psi(u-v) \mathrm{d} u \mathrm{~d} v \\
& +4 \operatorname{Tr} \int_{s}^{t} \int_{s}^{t} \mathrm{e}^{A u} \Phi \Phi^{\prime} \mathrm{e}^{A^{\prime} v} \psi(u-v) \mathrm{d} u \mathrm{~d} v \\
\leqslant & 4|A| \mathrm{e}^{2|A|} \sup _{s \in[n, n+1]}\left|\operatorname{Tr} Q_{s}\right|(t-s)^{2 H}+4 \operatorname{Tr} \sup _{n \leqslant u, v \leqslant n+1}\left|\mathrm{e}^{A u} \Phi \Phi^{\prime} \mathrm{e}^{A^{\prime} v}\right|(t-s)^{2 H} \\
\leqslant & K_{2}(t-s)^{2 H} .
\end{aligned}
$$

It follows from the Kolmogorov-Tchentsov Theorem that there exists a continuous modification $X_{t}^{n}$ of the process $X_{t}$ in the interval $[n, n+1]$. Therefore we can find sets $U_{n}, n \in \mathbb{N}$, such that

$$
X_{n+1}^{n}(\omega)=X_{n+1}(\omega)=X_{n+1}^{n+1}(\omega), \quad \omega \in \Omega \backslash U_{n}, \quad P\left(U_{n}\right)=0
$$

It follows that $P\left(\bigcup_{n=0}^{\infty} U_{n}\right)=0$ and for $\omega \in \Omega \backslash \bigcup_{n=0}^{\infty} U_{n}$ the process $\tilde{X}_{t}(\omega):=X_{t}^{n}(\omega)$ is a continuous modification of $X_{t}$ on $[n, n+1]$ (we will write $X_{t}$ for this modification).

Now we verify that the process (14) satisfies equation (13). We use the equation

$$
\begin{equation*}
\int_{0}^{t}\left(\int_{0}^{s} A \mathrm{e}^{A(s-r)} \Phi \mathrm{d} B^{H}(r)\right) \mathrm{d} s=\int_{0}^{t}\left(\int_{r}^{t} A \mathrm{e}^{A(s-r)} \mathrm{d} s\right) \Phi \mathrm{d} B^{H}(r) \tag{16}
\end{equation*}
$$

which we have got by integration by parts:

$$
\begin{aligned}
\int_{0}^{s} A \mathrm{e}^{A(s-r)} \Phi \mathrm{d} B^{H}(r) & =A \Phi B^{H}(s)+\int_{0}^{s} A^{2} \mathrm{e}^{A(s-r)} \Phi B^{H}(r) \mathrm{d} r \\
\int_{0}^{t}\left(\int_{r}^{t} A \mathrm{e}^{A(s-r)} \mathrm{d} s\right) \Phi \mathrm{d} B^{H}(r) & =\int_{0}^{t}\left(A+\int_{r}^{t} A^{2} \mathrm{e}^{A(s-r)} \mathrm{d} s\right) \Phi B^{H}(r) \mathrm{d} r .
\end{aligned}
$$

Using (16) we have for the left-side of (13)

$$
\begin{array}{rl}
x_{0}+\int_{0}^{t} A & A(s) \mathrm{d} s+\Phi B^{H}(t) \\
& =x_{0}+\int_{0}^{t} A\left(\mathrm{e}^{A s} x_{0}+\int_{0}^{s} \mathrm{e}^{A(s-r)} \Phi \mathrm{d} B^{H}(r)\right) \mathrm{d} s+\Phi B^{H}(t) \\
& =\mathrm{e}^{A t} x_{0}+\int_{0}^{t}\left(\int_{r}^{t} A \mathrm{e}^{A(s-r)} \mathrm{d} s\right) \Phi \mathrm{d} B^{H}(r)+\Phi B^{H}(t) \\
& =\mathrm{e}^{A t} x_{0}+\int_{0}^{t}\left(\mathrm{e}^{A(t-r)}-I\right) \Phi \mathrm{d} B^{H}(r)+\Phi B^{H}(t)=X(t)
\end{array}
$$

and hence $X_{t}$ is a strong solution. Let $Y_{t}$ be any strong solution, then

$$
\left|X_{t}-Y_{t}\right|=\left|\int_{0}^{t} A\left(X_{s}-Y_{s}\right) \mathrm{d} s\right| \leqslant \int_{0}^{t}|A|\left|X_{s}-Y_{s}\right| \mathrm{d} s \quad \text { for } t \geqslant 0 \quad \mathscr{P} \text {-a.s. }
$$

and using the Gronwall lemma we have $P\left[\left|X_{t}-Y_{t}\right|=0, t \geqslant 0\right]=1$.

Proposition 3.2. Let matrices $S(t), Q_{t} \in \mathbb{R}^{n \times n}$ be differentiable, $Q_{t}$ symmetric and positive definite, $S(t)$ invertible and let

$$
\begin{equation*}
p(t, x, y)=\frac{1}{(2 \pi)^{n / 2}\left(\operatorname{det} Q_{t}\right)^{1 / 2}} \exp \left\{-\frac{1}{2}(S(t) x-y)^{\prime} Q_{t}^{-1}(S(t) x-y)\right\} \tag{17}
\end{equation*}
$$

be the density of the $n$-dimensional Gaussian distribution $N_{n}\left(S(t) x, Q_{t}\right)$. Then

$$
\begin{align*}
\frac{\partial p(t, x, y)}{\partial t}= & \left\langle S(t)^{-1} \dot{S}(t) x, D_{x} p(t, x, y)\right\rangle  \tag{18}\\
& +\frac{1}{2} \operatorname{Tr}\left(S(t)^{-1} \dot{Q}_{t}\left(S(t)^{-1}\right)^{\prime} D_{x}^{2} p(t, x, y)\right)
\end{align*}
$$

Proof. We are using the following formulae for the time derivatives of the determinant and the inverse matrix:

$$
\frac{\mathrm{d}}{\mathrm{~d} t} \operatorname{det} Q_{t}=\operatorname{det} Q_{t} \operatorname{Tr}\left(Q_{t}^{-1} \dot{Q}_{t}\right), \quad \frac{\mathrm{d}}{\mathrm{~d} t} Q_{t}^{-1}=-Q_{t}^{-1} \dot{Q}_{t} Q_{t}^{-1}
$$

Thus for the first factor in the formula (17) we have

$$
\frac{\mathrm{d}}{\mathrm{~d} t}\left(\operatorname{det} Q_{t}\right)^{-1 / 2}=-\frac{1}{2} \frac{\operatorname{det} Q_{t} \operatorname{Tr} Q_{t}^{-1} \dot{Q}_{t}}{\left(\operatorname{det} Q_{t}\right)^{3 / 2}}=-\frac{1}{2} \frac{\operatorname{Tr} Q_{t}^{-1} \dot{Q}_{t}}{\left(\operatorname{det} Q_{t}\right)^{1 / 2}}
$$

and setting $z:=S(t) x-y$ we have

$$
\begin{aligned}
\frac{\partial p(t, x, y)}{\partial t}= & -p(t, x, y)\left\{\frac{1}{2} \operatorname{Tr}\left(Q_{t}^{-1} \dot{Q}_{t}\right)+(\dot{S}(t) x)^{\prime} Q_{t}^{-1} z\right\} \\
& +\frac{1}{2} p(t, x, y)\left\{z^{\prime} Q_{t}^{-1} \dot{Q}_{t} Q_{t}^{-1} z\right\}
\end{aligned}
$$

We have

$$
\begin{aligned}
& D_{x} p(t, x, y)=-p(t, x, y) S(t)^{\prime} Q_{t}^{-1} z \\
& D_{x}^{2} p(t, x, y)=p(t, x, y) S(t)^{\prime} Q_{t}^{-1} z z^{\prime} Q_{t}^{-1} S(t)-p(t, x, y) S(t)^{\prime} Q_{t}^{-1} S(t)
\end{aligned}
$$

and therefore

$$
\begin{aligned}
\left\langle S(t)^{-1} \dot{S}(t) x, D_{x} p(t, x, y)\right\rangle & =-p(t, x, y)(\dot{S}(t) x)^{\prime} Q_{t}^{-1} z \\
\left(S(t)^{-1}\right)^{\prime} D_{x}^{2} p(t, x, y) S(t)^{-1} \dot{Q}_{t} & =p(t, x, y) Q_{t}^{-1} z z^{\prime} Q_{t}^{-1} \dot{Q}_{t}-p(t, x, y) Q_{t}^{-1} \dot{Q}_{t}
\end{aligned}
$$

and

$$
\begin{aligned}
& \frac{1}{2} \operatorname{Tr}\left(\left(S(t)^{-1}\right)^{\prime} D_{x}^{2} p(t, x, y) S(t)^{-1} \dot{Q}_{t}\right) \\
& \quad=\frac{1}{2} p(t, x, y) z^{\prime} Q_{t}^{-1} \dot{Q}_{t} Q_{t}^{-1} z-\frac{1}{2} p(t, x, y) \operatorname{Tr}\left(Q_{t}^{-1} \dot{Q}_{t}\right) .
\end{aligned}
$$

Now it is easy to see that equation (18) is satisfied.
Let $X_{t}^{x}$ denote the solution of equation (12) with an initial condition $X_{0}^{x}=x$.

Proposition 3.3. Suppose $\Phi \Phi^{\prime}>0$ and let $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ be a Borel-measurable function satisfying the condition

$$
\begin{equation*}
\int_{-\infty}^{+\infty}|f(y)| \mathrm{e}^{-a|y|^{2}} \mathrm{~d} y<\infty \tag{19}
\end{equation*}
$$

for some $a>0$. Then there exists $t_{0}>0$ such that the function

$$
\begin{equation*}
u(t, x):=\mathbb{E} f\left(X_{t}^{x}\right) \tag{20}
\end{equation*}
$$

is well defined, differentiable with respect to $t$ for $t \in\left(0, t_{0}\right), u(t, \cdot) \in C^{\infty}\left(\mathbb{R}^{n}\right)$, and that

$$
\begin{equation*}
\frac{\partial u(t, x)}{\partial t}=\left\langle A x, D_{x} u(t, x)\right\rangle+\operatorname{Tr}\left(\Phi \Phi^{\prime} C(t) D_{x}^{2} u(t, x)\right) \tag{21}
\end{equation*}
$$

holds for $(t, x) \in\left(0, t_{0}\right) \times \mathbb{R}^{n}$ where $C(t)=\int_{0}^{t} \mathrm{e}^{-A^{\prime} s} \psi(s) \mathrm{d} s$ and $\psi(s)=H(2 H-$ $1) \mid s^{2 H-2}$. If $f$ is bounded and continuous, then

$$
\begin{equation*}
f(x)=\lim _{t \rightarrow 0+} u(t, x) \quad \forall x \in \mathbb{R}^{n} \tag{22}
\end{equation*}
$$

Proof. Let $p(t, x, y)$ be given by (17), where we substitute $S(t)=\mathrm{e}^{A t}$ and $Q_{t}$ according to (15). Note that

$$
S(t)^{-1} \dot{S}(t)=\mathrm{e}^{-A t} A \mathrm{e}^{A t}=A
$$

and due to the symmetry of the covariance matrix we have

$$
\begin{align*}
Q_{t} & =\int_{0}^{t} \int_{0}^{t} \mathrm{e}^{A u} \Phi \Phi^{\prime} \mathrm{e}^{A^{\prime} v} \psi(u-v) \mathrm{d} u \mathrm{~d} v \\
& =2 \int_{0}^{t} \mathrm{e}^{A u} \int_{0}^{u} \Phi \Phi^{\prime} \mathrm{e}^{A^{\prime} v} \psi(u-v) \mathrm{d} v \mathrm{~d} u \\
\dot{Q}_{t} & =2 \mathrm{e}^{A t} \Phi \Phi^{\prime} \int_{0}^{t} \mathrm{e}^{A^{\prime} v} \psi(t-v) \mathrm{d} v  \tag{23}\\
S(t)^{-1} \dot{Q}_{t}\left(S(t)^{-1}\right)^{\prime} & =2 \Phi \Phi^{\prime} \int_{0}^{t} \mathrm{e}^{A^{\prime}(v-t)} \psi(t-v) \mathrm{d} v  \tag{24}\\
& =2 \Phi \Phi^{\prime} \int_{0}^{t} \mathrm{e}^{-A^{\prime} u} \psi(u) \mathrm{d} u
\end{align*}
$$

Therefore, we have

$$
\begin{equation*}
\frac{\partial p(t, x, y)}{\partial t}=\left\langle A x, D_{x} p(t, x, y)\right\rangle+\operatorname{Tr}\left(\Phi \Phi^{\prime} C(t) D_{x}^{2} p(t, x, y)\right) \tag{25}
\end{equation*}
$$

Now we will prove that the function $u(t, x)=\int_{\mathbb{R}^{n}} f(y) p(t, x, y) \mathrm{d} y$ is differentiable. We bound eigenvalues $Q_{t}$ from above and from below to find an integrable majorant to derivatives of $p(t, x, y)$. Matrices $Q_{t}$ are positive definite for $t>0$, therefore there exist orthonormal matrices $U_{t}$ and diagonal matrices $\Lambda_{t}>0$ such that $Q_{t}=U_{t}^{\prime} \Lambda_{t} U_{t}$. If we denote the diagonal entries of the matrix $\Lambda_{t}$ by $\lambda_{i}(t)$, then $\lim _{t \rightarrow 0} \lambda_{i}(t)=0$ for $i \in\{1, \ldots, n\}$. It follows that we can find $t_{0}>0$ such that

$$
\begin{equation*}
0<\lambda_{i}(t)<(4 a)^{-1} \quad \text { for } 0<t<t_{0}, \quad i \in\{1, \ldots, n\} \tag{26}
\end{equation*}
$$

The function $\operatorname{det} Q_{t}$ is continuous, positive and so for every $t_{1} \in\left(0, t_{0}\right)$

$$
\begin{equation*}
\inf _{t_{1} \leqslant t \leqslant t_{0}} \operatorname{det} Q_{t}=: \inf _{t_{1} \leqslant t \leqslant t_{0}} \prod_{i=1}^{n} \lambda_{i}(t)=: \delta_{1}>0 \tag{27}
\end{equation*}
$$

and for $t \in\left(t_{1}, t_{0}\right)$ and for all $i$

$$
\lambda_{i}(t) \geqslant \delta_{1}\left(\prod_{j \neq i} \lambda_{j}(t)\right)^{-1}>\delta_{1}(4 a)^{n-1}=: \delta>0
$$

Therefore $\left|Q_{t}^{-1}\right|^{2}=\sum_{i=1}^{n} \lambda_{i}(t)^{-2}<n \delta^{-2}$ and we can estimate the normalizing constant in $p(t, x, y)$ by $(2 \pi)^{-n / 2}\left(\operatorname{det} Q_{t}\right)^{-1 / 2}<(2 \pi)^{-n / 2} \delta_{1}^{-1 / 2}$ for $t \geqslant t_{1}$. Further we estimate the exponent in $p(t, x, y)$. Let $|x|<k$ for some constant $k$. Then for $0<t<t_{0}$

$$
\begin{align*}
-\frac{1}{2}\left(\mathrm{e}^{A t} x-y\right)^{\prime} Q_{t}^{-1}\left(\mathrm{e}^{A t} x-y\right) & =-\frac{1}{2}\left(\mathrm{e}^{A t} x-y\right)^{\prime} U_{t}^{\prime} \Lambda_{t}^{-1} U_{t}\left(\mathrm{e}^{A t} x-y\right)  \tag{28}\\
& =-\frac{1}{2}\left(U_{t} \mathrm{e}^{A t} x-U_{t} y\right)^{\prime} \Lambda_{t}^{-1}\left(U_{t} \mathrm{e}^{A t} x-U_{t} y\right) \\
& \leqslant 4 a\left(U_{t} \mathrm{e}^{A t}\right)^{\prime} U_{t} y-2 a\left(U_{t} y\right)^{\prime} U_{t} y \\
& \leqslant 4 a\left|U_{t}\right|^{2} \mathrm{e}^{|A| t}|x||y|-2 a|y|^{2} \\
& \leqslant 4 a n \mathrm{e}^{|A| t_{0}} k|y|-2 a|y|^{2}  \tag{29}\\
& =-2 a|y|^{2}\left(1-\frac{2 n k \mathrm{e}^{|A| t_{0}}}{|y|}\right) \leqslant-a|y|^{2}
\end{align*}
$$

for $|y| \geqslant y_{0}, y_{0}$ sufficiently large. If we take $J=\left(t_{1}, t_{0}\right) \times B_{k}(0)\left(B_{k}(0)\right.$ is the open ball with center 0 and radius $k$ ), then we have

$$
\begin{align*}
& \left|\frac{\partial f(y) p(t, x, y)}{\partial t}\right| \vee\left|\frac{\partial^{k+l} f(y) p(t, x, y)}{\partial x_{i}^{k} \partial x_{j}^{l}}\right| \leqslant K|f(y)| \mathrm{e}^{-a|y|^{2}}  \tag{30}\\
& \forall y \in \mathbb{R}^{n}, \quad \forall(t, x) \in J, \quad i, j \in\{1, \ldots, n\}, \quad k, l \in\{0,1\}
\end{align*}
$$

where $K$ is a constant independent of $t$ and $x$. Our statement follows from (25) and the theorem about differentiation of an integral with respect to a parameter.

If the function $f$ is bounded and continuous then for almost all $\omega$ we have $f\left(X_{t}^{x}(\omega)\right) \rightarrow f(x)<\infty$ for all $x$ and the equation (22) follows from the dominated convergence theorem.

Corollary 3.4. If $\Phi \Phi^{\prime}>0, f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is a Borel-measurable function and there exists $b>0$ such that

$$
\begin{equation*}
\int_{-\infty}^{+\infty}|f(y)| \mathrm{e}^{-b|y|} \mathrm{d} y<\infty \tag{31}
\end{equation*}
$$

then the statement of Proposition 3.4 is valid for all $t_{0}>0$.

Proof. For arbitrary $t_{0}>0$ we have $\sup _{t \in\left(0, t_{0}\right)} \lambda_{i}(t)<\infty$. Therefore we can find $a>0$ which satisfies (26) and also (19) because

$$
\int_{-\infty}^{+\infty}|f(y)| \mathrm{e}^{-a|y|^{2}} \mathrm{~d} y=\int_{-\infty}^{+\infty} \mathrm{e}^{-a|y|^{2}+b|y|}|f(y)| \mathrm{e}^{-b|y|} \mathrm{d} y<\infty
$$

for every $a>0$.

## 4. Limit measure for the solution of the equation

In this section we are interested in weak and strong convergence of the probability laws of the solutions to the equation (12). In the case when the matrix $A$ is negatively definite there exists a limit measure in the weak sense. Let $N_{n}(a, Q)$ denote the $n$-dimensional Gaussian distribution with mean $a$ and a covariance matrix $Q$ and let $\mu_{t}^{x}$ denote the distribution of the solution $X_{t}^{x}$ to the equation (12) with an initial condition $X_{0}^{x}=x$.

Proposition 4.1. If $A$ is negatively definite then the matrix

$$
Q_{\infty}=\int_{0}^{\infty} \int_{0}^{\infty} \mathrm{e}^{A u} \Phi \Phi^{\prime} \mathrm{e}^{A^{\prime} v} \psi(u-v) \mathrm{d} u \mathrm{~d} v
$$

is well defined and for $\mu_{\infty}:=N\left(0, Q_{\infty}\right)$ we have

$$
\lim _{t \rightarrow \infty} \mu_{t}^{x}=\mu_{\infty}, \quad w^{\star} \text {-weakly }
$$

for every $x$.
Proof. If $A$ is negatively definite then there exist constants $M, \omega>0$ such that

$$
\left|\mathrm{e}^{A t}\right| \leqslant M \mathrm{e}^{-\omega t}, \quad t \geqslant 0
$$

We define a random process $Z(t)$ by

$$
Z(t):=\int_{0}^{t} \mathrm{e}^{A u} \Phi \mathrm{~d} B^{H}(u)
$$

which has the distribution $\mu_{t}=N\left(0, Q_{t}\right)$, where $Q_{t}$ is the covariance matrix from Proposition 3.2. Now we show that $(Z(t), t \geqslant 0)$ converge in $L^{2}\left(\Omega ; \mathbb{R}^{n}\right)$ for $t \rightarrow \infty$.

Let $t>s \geqslant 1$. There exist constants (all of them are denoted by the same $K$ ) such that

$$
\begin{aligned}
\mathbb{E}|Z(t)-Z(s)|^{2} & =\mathbb{E}\left|\int_{s}^{t} \mathrm{e}^{A u} \Phi \mathrm{~d} B^{H}(u)\right|^{2} \\
& =\operatorname{Tr} \int_{s}^{t} \int_{s}^{t} \mathrm{e}^{A u} \Phi \Phi^{\prime} \mathrm{e}^{A^{\prime} v} \psi(u-v) \mathrm{d} u \mathrm{~d} v \\
& \leqslant \int_{s}^{t} \int_{s}^{t}\left|\mathrm{e}^{A u} \Phi\right|\left|\mathrm{e}^{A^{\prime} v} \Phi\right| \psi(u-v) \mathrm{d} u \mathrm{~d} v \\
& \leqslant K \int_{s}^{\infty} \int_{s}^{\infty} \mathrm{e}^{-\omega u} \mathrm{e}^{-\omega v} \psi(u-v) \mathrm{d} u \mathrm{~d} v \\
& =2 K \int_{s}^{\infty} \int_{s}^{v} \mathrm{e}^{-\omega u} \mathrm{e}^{-\omega v} \psi(u-v) \mathrm{d} u \mathrm{~d} v \\
& \leqslant K \int_{s}^{\infty} \mathrm{e}^{-\omega v} \int_{s}^{v} \mathrm{e}^{-\omega u}(v-u)^{2 H-2} \mathrm{~d} u \mathrm{~d} v \\
& \leqslant K \int_{s}^{\infty} \mathrm{e}^{-\omega v} v^{2 H-1} \int_{0}^{1}(1-u)^{2 H-2} \mathrm{~d} u \mathrm{~d} v \\
& \leqslant K \int_{s}^{\infty} \mathrm{e}^{-\omega v} v^{2 H-1} \mathrm{~d} v
\end{aligned}
$$

The last integral on the right-hand side converges to 0 for $s \rightarrow \infty$. Therefore every sequence $Z(t), t \rightarrow \infty$ is Cauchy sequence in $L^{2}\left(\Omega ; \mathbb{R}^{n}\right)$ and hence there exists a random vector $Z(\infty)$ such that $Z(t) \rightarrow Z(\infty)$ in $L^{2}(\Omega)$. The probability law of $Z(\infty)$ is $\mu_{\infty}=N\left(0, Q_{\infty}\right)$, because $Q_{t} \rightarrow Q_{\infty}$. For $x \in \mathbb{R}^{n}$ and $\varphi: \mathbb{R}^{n} \rightarrow \mathbb{R}$ a bounded and lipschitz continuous function we have

$$
\begin{aligned}
\left|\int_{\mathbb{R}^{n}} \varphi \mathrm{~d} \mu_{t}^{x}-\int_{\mathbb{B}^{n}} \varphi \mathrm{~d} \mu_{\infty}\right| & =\left|\mathbb{E} \varphi\left(\mathrm{e}^{A t} x+Z(t)\right)-\mathbb{E} \varphi(Z(\infty))\right| \\
& \leqslant k_{\varphi}\left(\left|\mathrm{e}^{A t} x\right|+\mathbb{E}|Z(t)-Z(\infty)|\right) \\
& \left.\leqslant k_{\varphi}\left(M|x| \mathrm{e}^{-\omega t}+\mathbb{E}|Z(t)-Z(\infty)|^{2}\right)^{1 / 2}\right)
\end{aligned}
$$

where $k_{\varphi}$ is the Lipschitz constant of $\varphi$. Our assertion easily follows.
In what follows we consider the metric of total variation $\left\|P_{1}-P_{2}\right\|_{\text {var }}:=$ $\sup _{B \in \mathcal{B}}\left|P_{1}(B)-P_{2}(B)\right|$ as the "distance" between two probability measures $P_{1}$ and $P_{2}$ on $\mathcal{B} \equiv \mathcal{B}\left(\mathbb{R}^{n}\right)$. If these measures are absolutely continuous with respect to the Lebesgue measure with densities $f_{1}$ and $f_{2}$ we have

$$
\begin{equation*}
\left\|P_{1}-P_{2}\right\|_{\mathrm{var}}=\frac{1}{2} \int_{\mathbb{R}^{n}}\left|f_{1}(x)-f_{2}(x)\right| \mathrm{d} x . \tag{32}
\end{equation*}
$$

Suppose moreover that $A=A^{\prime}$ and $\Phi \Phi^{\prime} A=A \Phi \Phi^{\prime}$. Then there exists an orthonormal basis ( $e_{i}$ ) and numbers $-\alpha_{n}<-\alpha_{n-1}<\ldots<-\alpha_{1}<0$ and $0<\lambda_{i}$ such that

$$
\begin{equation*}
A e_{i}=-\alpha_{i} e_{i} \quad \text { and } \quad \Phi \Phi^{\prime} e_{i}=\lambda_{i} e_{i} \tag{33}
\end{equation*}
$$

(see [9]). In this case the eigenvalues of $Q_{t}$ are

$$
\begin{align*}
q_{i}(t) & =\lambda_{i} \int_{0}^{t} \int_{0}^{t} \mathrm{e}^{-\alpha_{i} u} \mathrm{e}^{-\alpha_{i} v} \psi(u-v) \mathrm{d} u \mathrm{~d} v  \tag{34}\\
& =\frac{\lambda_{i}}{\alpha_{i}^{2 H}} \int_{0}^{t \alpha_{i}} \int_{0}^{t \alpha_{i}} \mathrm{e}^{-u} \mathrm{e}^{-v} \psi(u-v) \mathrm{d} u \mathrm{~d} v
\end{align*}
$$

and the eigenvalues of $Q_{\infty}$ are

$$
\begin{equation*}
q_{i}=\frac{\lambda_{i}}{\alpha_{i}^{2 H}} \int_{0}^{\infty} \int_{0}^{\infty} \mathrm{e}^{-u} \mathrm{e}^{-v} \psi(u-v) \mathrm{d} u \mathrm{~d} v \tag{35}
\end{equation*}
$$

If we denote

$$
f(y, z, t)=(\pi)^{-n / 2}\left(\operatorname{det} Q_{t}\right)^{-1 / 2} \exp \left\{-(z-y)^{\prime} Q_{t}^{-1}(z-y)\right\}
$$

for $y, z \in \mathbb{R}^{n}, t \in \mathbb{R}_{+} \cup\{\infty\}$ then the measure $\mu_{t}^{x}$ has the density $f\left(\cdot, \mathrm{e}^{A t} x, t\right)$ and the measure $\mu_{\infty}$ has the density $f(\cdot, 0, \infty)$.

Proposition 4.2. If $A$ and $\Phi \Phi^{\prime}$ satisfy the above conditions then

$$
\begin{equation*}
\left\|\mu_{t}^{x}-\mu_{\infty}\right\|_{\mathrm{var}} \leqslant K(|x|+1) \mathrm{e}^{-t \alpha}, \quad t \geqslant 1 \tag{36}
\end{equation*}
$$

for an arbitrary $0<\alpha<\alpha_{1}$ and a constant $K=K\left(H, \alpha, \lambda_{i}\right)$.
Proof. The measures $\mu_{t}^{x}$ and $\mu_{\infty}$ are absolutely continuous with respect to the Lebesgue measure, therefore

$$
\begin{align*}
\left\|\mu_{t}^{x}-\mu_{\infty}\right\|_{\text {var }}= & \frac{1}{2} \int_{\mathbb{R}^{n}}\left|f\left(y, \mathrm{e}^{A t} x, t\right)-f(y, 0, \infty)\right| \mathrm{d} y  \tag{37}\\
\leqslant & \int_{\mathbb{R}^{n}}\left|f\left(y, \mathrm{e}^{A t} x, t\right)-f(y, 0, t)\right| \mathrm{d} y \\
& +\int_{\mathbb{R}^{n}}\left|f(y, 0, t)-\left(\operatorname{det} Q_{t}^{-1} \operatorname{det} Q_{\infty}\right)^{1 / 2} f(y, 0, \infty)\right| \mathrm{d} y \\
& +\int_{\mathbb{R}^{n}}\left|\left(\operatorname{det} Q_{t}^{-1} \operatorname{det} Q_{\infty}\right)^{1 / 2} f(y, 0, \infty)-f(y, 0, \infty)\right| \mathrm{d} y .
\end{align*}
$$

Due to the mean value theorem and to the Fubini theorem we can estimate the first term on the right-hand side:

$$
\begin{align*}
\int_{\mathbb{R}^{n}} \mid f\left(y, \mathrm{e}^{A t} x,\right. & t)-f(y, 0, t) \mid \mathrm{d} y  \tag{38}\\
= & \int_{\mathbb{R}^{n}}\left|\int_{0}^{1}\left\langle D_{z} f\left(y, u \mathrm{e}^{A t} x, t\right), \mathrm{e}^{A t} x\right\rangle \mathrm{d} u\right| \mathrm{d} y \\
& =\int_{\mathbb{R}^{n}}\left|\int_{0}^{1}\left\langle Q_{t}^{-1}\left(u \mathrm{e}^{A t} x-y\right), \mathrm{e}^{A t} x\right\rangle f\left(y, u \mathrm{e}^{A t} x, t\right) \mathrm{d} u\right| \mathrm{d} y \\
& \leqslant\left|Q_{t}^{-1}\right|\left|\mathrm{e}^{A t} x\right| \int_{0}^{1} \int_{\mathbb{R}^{n}}\left|u \mathrm{e}^{A t} x-y\right| f\left(y, u \mathrm{e}^{A t} x, t\right) \mathrm{d} y \mathrm{~d} u
\end{align*}
$$

Now we substitute $z=u \mathrm{e}^{A t} x-y$ and use the Hölder inequality, which yields

$$
\begin{align*}
\left|Q_{t}^{-1}\right|\left|\mathrm{e}^{A t} x\right| \int_{\mathbb{R}^{n}}|z| f(z, 0, t) \mathrm{d} z & \leqslant\left|Q_{t}^{-1}\right|\left|Q_{t}\right|\left|\mathrm{e}^{A t} x\right|  \tag{39}\\
& \leqslant\left|Q_{1}^{-1}\right|\left|Q_{\infty}\right||x| \mathrm{e}^{-t \alpha_{1}}
\end{align*}
$$

The second term on the right-hand side in (37) we estimate as follows:

$$
\begin{align*}
& \int_{\mathbb{R}^{n}}\left|f(y, 0, t)-\left(\operatorname{det} Q_{t}^{-1} \operatorname{det} Q_{\infty}\right)^{1 / 2} f(y, 0, \infty)\right| \mathrm{d} y  \tag{40}\\
& =\int_{\mathbb{R}^{n}}\left(\pi^{n} \operatorname{det} Q_{t}\right)^{-1 / 2}\left|\exp \left\{-\sum_{i=1}^{n} q_{i}(t)^{-1} y_{i}^{2}\right\}-\exp \left\{-\sum_{i=1}^{n} q_{i}^{-1} y_{i}^{2}\right\}\right| \mathrm{d} y \\
& =\left(\pi^{n} \operatorname{det} Q_{t}\right)^{-1 / 2} \int_{\mathbb{R}^{n}}\left|\int_{0}^{1} \frac{\partial}{\partial u} \exp \left\{-\sum_{i=1}^{n}\left(q_{i}^{-1}+u\left(q_{i}(t)^{-1}-q_{i}^{-1}\right)\right) y_{i}^{2}\right\} \mathrm{d} u\right| \mathrm{d} y \\
& =\left(\pi^{n} \operatorname{det} Q_{t}\right)^{-1 / 2} \sum_{i=1}^{n} \int_{\mathbb{R}^{n}}\left(\frac{q_{i}-q_{i}(t)}{q_{i}(t) q_{i}}\right) y_{i}^{2} \exp \left\{-\sum_{i=1}^{n} q_{i}^{-1} y_{i}^{2}\right\} \\
& \times \int_{0}^{1} \exp \left\{-u \sum_{i=1}^{n}\left(q_{i}(t)^{-1}-q_{i}^{-1}\right) y_{i}^{2}\right\} \mathrm{d} u \mathrm{~d} y \\
& \leqslant \\
& \leqslant \max _{i \in\{1, \ldots, n\}}\left\{\frac{q_{i}-q_{i}(t)}{q_{i} q_{i}(t)}\right\}\left(\pi^{n} \operatorname{det} Q_{t}\right)^{-1 / 2} \int_{\mathbb{R}^{n}}|y|^{2} \exp \left\{-\sum_{i=1}^{n} q_{i}^{-1} y_{i}^{2}\right\} \mathrm{d} y \\
& =\max _{i \in\{1, \ldots, n\}}\left\{\frac{q_{i}-q_{i}(t)}{q_{i} q_{i}(t)}\right\}\left(\operatorname{det} Q_{t}^{-1} Q_{\infty}\right)^{1 / 2}\left|Q_{\infty}\right| \\
& \leqslant \sum_{i=1}^{n}\left(\frac{q_{i}-q_{i}(t)}{q_{i} q_{i}(t)}\right)\left(\operatorname{det} Q_{1}^{-1} Q_{\infty}\right)^{1 / 2}\left|Q_{\infty}\right|,
\end{align*}
$$

because $0<q(1)<q_{i}(t)$ for all $i$ and for all $t \geqslant 1$. Let $\alpha<\alpha_{1}$, then there exist constants $k_{i}$ (their values may differ from line to line) depending only on $H, \lambda_{i}, \alpha_{i}$
and $\alpha$ such that

$$
\begin{align*}
q_{i}-q_{i}(t) & =k_{i} \int_{t}^{\infty} \mathrm{e}^{-\alpha_{i} u} \int_{0}^{u} \mathrm{e}^{-\alpha_{i} v}|u-v|^{2 H-2} \mathrm{~d} v \mathrm{~d} u  \tag{41}\\
& \leqslant k_{i} \int_{t}^{\infty} \mathrm{e}^{-\alpha_{i} u} \int_{0}^{\infty} \mathrm{e}^{-\alpha_{i} v}|u-v|^{2 H-2} \mathrm{~d} v \mathrm{~d} u \\
& =k_{i} \int_{t}^{\infty} \mathrm{e}^{-\alpha_{i} u} \int_{0}^{\infty} \mathrm{e}^{-\alpha_{i} v} u^{2 H-2}|1-v / u|^{2 H-2} \mathrm{~d} v \mathrm{~d} u \\
& =k_{i} \int_{t \alpha_{i}}^{\infty} \mathrm{e}^{-u} u^{2 H-1} \int_{0}^{\infty} \mathrm{e}^{-u v}|v-1|^{2 H-2} \mathrm{~d} v \mathrm{~d} u \\
& \leqslant k_{i} \int_{t \alpha_{i}}^{\infty} \mathrm{e}^{-u} u^{2 H-1} \mathrm{~d} u=k_{i} \int_{t}^{\infty} \mathrm{e}^{-\alpha_{i} u} u^{2 H-1} \mathrm{~d} u \\
& =k_{i} \mathrm{e}^{-t \alpha} \int_{t}^{\infty} \mathrm{e}^{-\left(\alpha_{i}-\alpha\right) u} u^{2 H-2} \mathrm{~d} u \leqslant k_{i} \mathrm{e}^{-t \alpha} .
\end{align*}
$$

Now we can estimate the expression (40) by a negative exponential, because

$$
\begin{equation*}
\sum_{i=1}^{n}\left(\frac{q_{i}-q_{i}(t)}{q_{i} q_{i}(t)}\right) \leqslant \mathrm{e}^{-t \alpha} \sum_{i=1}^{n} \frac{k_{i}}{q_{i} q_{i}(1)} . \tag{42}
\end{equation*}
$$

For the third term in (37) we have

$$
\begin{align*}
\int_{\mathbb{R}^{n}}\left|\left(\operatorname{det} Q_{t}^{-1} \operatorname{det} Q_{\infty}\right)^{1 / 2} f(y, 0, \infty)-f(y, 0, \infty)\right| \mathrm{d} y  \tag{43}\\
\quad=\left|\left(\operatorname{det} Q_{t}^{-1} \operatorname{det} Q_{\infty}\right)^{1 / 2}-1\right|=\left|\prod_{i=1}^{n} \sqrt{q_{i} q_{i}(t)^{-1}}-1\right|
\end{align*}
$$

As we have

$$
\begin{aligned}
\sqrt{\frac{q_{i}}{q_{i}(t)}}=\frac{\sqrt{q_{i}}-\sqrt{q_{i}(t)}}{\sqrt{q_{i}(t)}}+1 & =\frac{q_{i}-q_{i}(t)}{\sqrt{q_{i}(t)}\left(\sqrt{q_{i}}+\sqrt{q_{i}(t)}\right)}+1 \\
& \leqslant \frac{q_{i}-q_{i}(t)}{q_{i}(t)}+1 \leqslant \exp \left\{\frac{q_{i}-q_{i}(t)}{q_{i}(t)}\right\}
\end{aligned}
$$

and $0<q_{i}-q_{i}(t) \rightarrow 0$ for $t \rightarrow+\infty$ we get

$$
\begin{equation*}
\left|\prod_{i=1}^{n} \sqrt{q_{i} q_{i}(t)^{-1}}-1\right| \leqslant \exp \left\{\sum_{i=1}^{n} \frac{q_{i}-q_{i}(t)}{q_{i}(t)}\right\}-1 \leqslant K_{1} \sum_{i=1}^{n} \frac{q_{i}-q_{i}(t)}{q_{i}(t)} \tag{44}
\end{equation*}
$$

for some constant $K_{1}$. Similarly to (42) there exists a constant $K$ such that $K \mathrm{e}^{-t \alpha}$, $t \geqslant 1$ is a majorant for the right-hand side of (44).

Remark 4.3. In particular, if $A$ and $\Phi$ are diagonal, Proposition 4.2 may be applied.

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