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Fractional Brownian Motion and Applications in Financial Mathematics

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Name: Zhu Vorname: Ke

geboren am: 03.12.1985 in: Wuhan

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Betreuer: Prof. Dr. rer. nat. Martin Keller-Ressel

Abstract

Fractional Brownian motion (fBm) $(U_H(t))_{t\in\mathbb{R}}$ has an integration representation defined by Mandelbrot and Van Ness[17]

$$U_H(t) - U_H(s) = \frac{1}{\Gamma(H + \frac{1}{2})} \left(\int_{\mathbb{R}} \mathbb{1}_{\{t \ge u\}} \cdot (t - u)^{H - \frac{1}{2}} - \mathbb{1}_{\{s \ge u\}} (-u)^{H - \frac{1}{2}} dB_u \right),$$

where $(B_t)_{t \in \mathbb{R}}$ is two-sides Brownian motion.

One of applications of fBm is fractional Ornstein-Uhlenbeck process (fOU)

$$X_t = e^{-at} \left(\gamma \int_{-\infty}^t e^{au} dU_H(u) \right),$$

which is as the stationary solution of the SDE with

$$dX_t = -aX_t dt + \gamma dU_H(t),$$

where $a, \gamma \in \mathbb{R}_+$.

This thesis is devoted to the study of fBm, fOU and financial modelings, which are derived from fOU.

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1 Introduction

The aim of this Diploma thesis is to study fractional Brownian motion (fBm). Unlike ordinary Brownian motion, fBm does not need to have independent increments. fBm was first introduced by Kolmogorov[11] as a centered Gaussian process, which has covariance function as follows

$$Cov[U_H(t), U_H(s)] = \frac{1}{2}(|t|^{2H} + |s|^{2H} - |t - s|^{2H}),$$

where H is real number in (0,1). The successive pioneer work can be traced to Mandelbrot and Van Ness[17]. Where fBm $(U_H(t))_{t\in\mathbb{R}}$ has an integration representation

$$U_H(t) - U_H(s) = \frac{1}{\Gamma(H + \frac{1}{2})} \left(\int_{\mathbb{R}} \mathbb{1}_{\{t \ge u\}} \cdot (t - u)^{H - \frac{1}{2}} - \mathbb{1}_{\{s \ge u\}} (-u)^{H - \frac{1}{2}} dB_u \right),$$

and has stationary self-similar increments satisfying for $t \in \mathbb{R}, \tau > 0$

$$U_H(t+\tau) - U_H(t) \sim \kappa \tau^H$$

with some κ . According to Theorem of Kolmogorov Chentsov, fBm must be a continuous-time process.

One of applications of fBm is fractional Ornstein-Uhlenbeck process (fOU)

$$X_t = e^{-at} \left(\gamma \int_{-\infty}^t e^{au} \, dU_H(u) \right)$$

which is as the stationary solution of the SDE with

$$dX_t = -aX_t dt + \gamma dU_H(t)$$

where $a, \gamma \in \mathbb{R}_+$. Although fBm is not a semimartingale for $H \neq \frac{1}{2}$ (Liptser & Shiryayev[12]), Cheridito prove fOU as integral driven by $U_H(t)$ in sense of Riemann-Stieljes exists, i.e., for a > 0,

$$\int_0^t e^{au} dU_H(u)$$

is well-defined. In addition, if $H \in (\frac{1}{2}, 1)$, fOU exhibits long memory property. More recently, fOU is more and more commonly applied in stochastic volatility models. In particular, Comte and Renault[7] assume the log-volatility to follow fOU with $H \in (\frac{1}{2}, 1)$. In contrast, Gatheral et al.[20] take $H \in (0, \frac{1}{2})$.

The structure of this thesis is as follows. We provide in Section 2 basic concepts of probability theory and stochastic process involving Gaussian process which could be characterized

by characteristic function. Brownian motion is discussed at the end of this section as an example. Section 3 focuses on stable integrals, which ensure that fBm can be defined as a Gaussian process. Section 4 presents definition of fBm in sense of stable integral, based on which we show self-similarity, stationary property of increments of fBm, regularity, and nonsemimartingale (except for $H=\frac{1}{2}$). From Section 5 up to the end we turn towards applications of FBM in financial mathematics. In Section 5, we deal with fOU, which exhibits long memory if $H > \frac{1}{2}$. In Section 6 we list out applications of FBM, not only for pricing models of risky asset (fractional Black-Scholes) but also for stochastic volatility models. We complete the proof of Cheridito to show that if a minimal amount of time between two successive transactions exists, fractional Black-Scholes would be arbitragefree. In fractional stochastic volatility model, we use fOU to model log-volatility. It differs from the choise of H. Whilst $H > \frac{1}{2}$ (in FSV) can ensure long memory, on the other hand, if $H<\frac{1}{2}$ the model (RFSV) can generate more desirable volatility smoothness according to empirical data. In order to combine this two characteristics, we define weighted fractional stochatic volatility model which inherits long memory property from FSV and could have a very close result as in RFSV.

2 Gaussian Processes and Brownian Motion

In this section we start off by looking at some general concepts of probability spaces and stochastic processes, in which the Gaussian process is an important example. Within the framework of Gaussian processes, one could specify a stationary and independent behavior. This therefore leads us to introduction of the Brownian motion.

2.1 Probability Spaces and Stochastic Processes

DEFINITION 2.1. Let \mathscr{A} be a collection of subsets of a set Ω . \mathscr{A} is said to be a σ - Algebra on Ω , if it satisfies the following conditions:

- (i) $\Omega \in \mathscr{A}$.
- (ii) For any set $F \in \mathcal{A}$, its complement $F^c \in \mathcal{A}$.
- (iii) If there is a series $\{F_n\}_{n\in\mathbb{N}}$ such that $\{F_n\}_{n\in\mathbb{N}}\subseteq\mathscr{A}$, then $\cup_{n\in\mathbb{N}}F_n\in\mathscr{A}$.

DEFINITION 2.2. A mapping \mathcal{P} is said to be a *probability measure* from \mathscr{A} to $\mathscr{B}(\mathbb{R}^n)$, if $\mathcal{P}\left[\sum_{n=1}^{\infty}F_n\right]=\sum_{n=1}^{\infty}\mathcal{P}\left[F_n\right]$ for any $\{F_n\}_{n\in\mathbb{N}}$ disjoint in \mathscr{A} satisfying $\sum_{n=1}^{\infty}F_n\in\mathscr{A}$.

DEFINITION 2.3. A probability space is defined as a triple $(\Omega, \mathscr{A}, \mathcal{P})$ of a set Ω , a σ -Algebra \mathscr{A} of Ω and a measure \mathcal{P} from \mathscr{A} to $\mathscr{B}(\mathbb{R}^n)$.

The σ - Algebra generated of all open sets on \mathbb{R}^n is called the *Borel* σ - Algebra which is denoted by $\mathscr{B}(\mathbb{R}^n)$. Let μ be a probability measure on \mathbb{R}^n . Indeed, $(\mathbb{R}^n, \mathscr{B}(\mathbb{R}^n), \mu)$ is a special case that probability space on \mathbb{R}^n . A function f mapping from $(\mathcal{D}, \mathcal{D}, \mu)$ into $(\mathcal{E}, \mathcal{E}, \nu)$ is measurable, if its collection of the inverse image of \mathcal{E} is a subset of \mathcal{D} . A random variable is a \mathbb{R}^n -valued measurable function on some probability space. Let \mathcal{P} represent a probability measure, recall that in probability theory, for $B \in \mathscr{B}(\mathbb{R}^n)$ we call $\mathcal{P}[\{X \in B\}]$ the distribution of X. We write also $\mathcal{P}_X[\cdot]$ or $\mathcal{P}[X]$ for convenience for those notations.

DEFINITION 2.4. Let $(\Omega, \mathscr{A}, \mathcal{P})$ be a probability space. A *n*-dimensional stochastic process $(X_t)_t$ is a family of random variable such that $X_t(\omega): \Omega \longrightarrow \mathbb{R}^n, \forall t \in T$, where T denotes the set of Index of Time.

Without specification, we set $T = \mathbb{R}$. Some basic definitions, which are needed in following sections, are given.

DEFINITION 2.5. A stochastic process $(X_t)_{t\in T}$ is said to be *stationary*, if the joint distribution

$$\mathcal{P}\left[X_{t_1},\ldots,X_{t_n}\right] = \mathcal{P}\left[X_{t_1+\tau},\ldots,X_{t_n+\tau}\right]$$

for t_1, \ldots, t_n and $t_1 + \tau, \ldots, t_n + \tau \in \mathbb{R}$.

DEFINITION 2.6. Let $(X_t)_t$ be a stochastic process.

$$\varsigma_X(t,s) := \operatorname{Cov}(X_t, X_s)$$

is called autocovariance between s, t and

$$\eta_X(t,s) := \frac{\mathrm{Cov}[\mathbf{X}_{\mathrm{t}},\mathbf{X}_{\mathrm{s}}]}{\sqrt{\mathrm{Var}[X_t]\mathrm{Var}[X_s]}}$$

is called autocorrelation betwenn s, t.

DEFINITION 2.7. A stochastic process $(X_t)_t$ is said to be weak stationary if

$$E[X_t] = E[X_{t+\tau}]$$

and

$$\varsigma_X(t,s) = \varsigma_X(t-s,0)$$

for $\tau, s \in \mathbb{R}$.

Remark that, weak stationarity is more general than stationarity. If $(X_t)_t$ is a weak stationary process, for any t, we write $\varsigma_X(\tau)$ for $\varsigma_X(t+\tau,t)$. $\eta_X(\tau)$ is used in the same way.

We use a notation $X \sim Y$ represents X equals Y in distribution.

DEFINITION 2.8. A stochastic process $(X_t)_t$ is said to be α -self similar if $(X_{ct_1}, \ldots, X_{ct_k}) \sim (c^{\alpha}X_{t_1}, \ldots, c^{\alpha}X_{t_k})$ for any $t_1, \ldots, t_k \in \mathbb{R}$ and $c \in \mathbb{R}_+$.

2.2 Normal Distribution and Gaussian Process

DEFINITION 2.9 (1-dimensional normal distribution). A \mathbb{R} -valued random variable X is said to be *standard normal distributed* or *standard Gaussian*, if its distribution can be described as

$$\mathcal{P}[X \le x] = (2\pi)^{-\frac{1}{2}} \int_{-\infty}^{x} e^{-\frac{u^2}{2}} du$$
 (2.1)

for $x \in \mathbb{R}$.

The integrand of (2.1) is also called *density function* of a standard Gaussian random variable.

DEFINITION 2.10. A \mathbb{R} -valued random variable X is said to be normally distributed or Gaussian with a expected value μ and a variance σ^2 , if

$$(X-\mu)/\sigma$$

is standard Gaussian for $\sigma > 0$.

PROPOSITION 2.11. Let X be a \mathbb{R} -valued Gaussian random variable with expected value μ and variance σ^2 , then it is distributed as

$$\mathcal{P}[X \le x] = (2\pi\sigma^2)^{-\frac{1}{2}} \int_{-\infty}^{x} e^{-\frac{(u-\mu)^2}{2\sigma^2}} du$$

Proof. Suppose $X = \sigma Y + \mu$ with Y standard Gaussian. We denote this mapping by $g(y): y \to \sigma y + \mu$ and give the inverse $g^{-1}(x): x \to \frac{(x-\mu)}{\sigma}$. The distribution function of X is

$$\int_{\Omega} \mathcal{P}[X \in dx] = \int_{\Omega} \mathcal{P}[Y \circ g \in dx]$$

$$= \int_{\mathbb{R}^{\circ g}} f_Y \circ g^{-1}(x) dx$$

$$= \int_{\mathbb{R}} \sigma \frac{1}{\sqrt{2\pi}} \exp\{\frac{(\frac{(x-\mu)}{\sigma})^2}{2}\} dy$$

$$= (2\pi\sigma^2)^{-\frac{1}{2}} \int_{\mathbb{R}} \exp\{-\frac{(x-\mu)^2}{2\sigma^2}\} dy,$$

where f_Y is density function of Y.

It is denoted by $X \sim (2\pi)^{-\frac{1}{2}} e^{-\frac{x^2}{2}} dx$, if X is standard Gaussian. In order to verifying the behavior of a normal distributed random variable we use the characteristic function in probability theory, cf.[1].

THEOREM 2.12. Let X be a \mathbb{R} -valued Gaussian random variable with expected value μ and variance σ^2 . The characteristic function of X

$$\Psi_X(\xi) := \int_{\mathbb{R}} e^{ix\xi} \mathcal{P}[X \in dx] = e^{i\mu\xi - \frac{1}{2}(\sigma\xi)^2}$$
 (2.2)

for $\xi \in \mathbb{R}$.

Proof. Cf.[19]. We assume firstly Y is standard Gaussian. In terms of the Definion of characteristic function of a standard Gaussian Y, integrating its density function over \mathbb{R} we get

$$\Psi_Y(\xi) = \int_{\mathbb{R}} (2\pi)^{-\frac{1}{2}} e^{-\frac{y^2}{2}} e^{iy\xi} dy,$$

take differentiating both sides of the equation by ξ , then

$$\Psi'_{Y}(\xi) = \int_{\mathbb{R}} (2\pi)^{-\frac{1}{2}} e^{-\frac{y^{2}}{2}} e^{iy\xi} ix \, dy$$

$$= (-i) \cdot \int_{\mathbb{R}} (2\pi)^{-\frac{1}{2}} (\frac{d}{dx} e^{-\frac{y^{2}}{2}}) e^{iy\xi} \, dy$$

$$\stackrel{part.int.}{=} - \int_{\mathbb{R}} (2\pi)^{-\frac{1}{2}} e^{-\frac{y^{2}}{2}} e^{iy\xi} \xi \, dy$$

$$= -\xi \Psi_{Y}(\xi).$$

for $\xi \in \mathbb{R}$. Obviously, $\Psi(\xi) = \Psi(0)e^{-\frac{\xi^2}{2}}$ is the solution of the partial differential equation, and $\Psi(0)$ equals 1, hence $\Psi(\xi) = e^{-\frac{\xi^2}{2}}$.

Let
$$X = \sigma Y + \mu$$

$$\Psi_X(\xi)$$

$$= \operatorname{E}[e^{i\xi X}]$$

$$= \operatorname{E}[e^{i\xi(\sigma Y + \mu)}]$$

$$= e^{i\xi\mu} \operatorname{E}[e^{iy\xi(\sigma Y)}]$$

$$= e^{i\xi\mu} \operatorname{E}[e^{iy(\xi\sigma)Y}]$$

$$= e^{i\xi\mu} \Psi_Y(\xi\sigma)$$

$$= e^{i\mu\xi - \frac{1}{2}(\sigma\xi)^2}$$

DEFINITION 2.13. Let X be a \mathbb{R}^n -valued random vector. X is said to be normally distributed or Gaussian, if for any $d \in \mathbb{R}^n$ such that $d^T X$ is Gaussian in \mathbb{R} .

DEFINITION 2.14. A stochastic process $(X_t)_{t\in T}$ is said to be *Gaussian process* if the joint distribution of any finite instance is Gaussian, that means $(X_{t_1}, \ldots, X_{t_n})$ has joint Gaussian distribution in \mathbb{R}^n for $t_1, \ldots, t_n \in T$.

The definition immediately shows every instances X_t in Gaussian process is Gaussian.

COROLLARY 2.15. Let $(X_t)_{t\in T}$ be a stochastic process. The following condition is equivalent to Definition 2.14.

$$\sum_{j}^{n} c_{t_j} X_{t_j}$$

is Gaussian for any $t_1, \ldots, t_n \in T, c_{t_j} \in \mathbb{R}$ for $j \in 1, \ldots, n$.

Proof. It is clear due to Definition 2.13.

LEMMA 2.16. Let X be a \mathbb{R}^n -valued normally distributed ramdon vector. Then its characteristic function is

$$E e^{i\xi^T X} = e^{i\xi^T m - \frac{1}{2}\xi^T \Sigma \xi}.$$
 (2.3)

For $\xi \in \mathbb{R}^n$. Where $m \in \mathbb{R}^n$, $\Sigma \in \mathbb{R}^{n \times n}$ are mean vector, covariance matrix of X respectively. Furthermore, the density function of X is

$$(2\pi)^{-\frac{n}{2}} \left(\det \Sigma\right)^{-\frac{1}{2}} e^{-\frac{1}{2}(x-m)^T \Sigma^{-1}(x-m)}.$$
 (2.4)

Remark, the equation (2.3) can also be as definition of characteristic function of a n-dimensional normally distributed random variable. I.e., any normally distributed random variable can be characterized by form of the equation (2.3).

Proof. Since X normally distributed on \mathbb{R}^n , then $\xi^T X$ is normally distributed on \mathbb{R} . Due to the Theorem 2.12, there is

$$\begin{split} \mathbf{E}e^{i\xi^TX} &= \mathbf{E}e^{i\cdot 1\cdot \xi^TX} \\ &= e^{i\mathbf{E}\left[\xi^TX\right] - \frac{1}{2}\mathrm{Var}\left[\xi^TX\right]} \\ &= e^{i\xi^T\mathbf{E}\left[X\right] - \frac{1}{2}\xi^T\mathrm{Var}\left[X\right]\xi} \\ &= e^{i\xi^Tm - \frac{1}{2}\xi^T\Sigma\xi}. \end{split}$$

Moreover, since Σ symmetric and positive definit, there exist $\Sigma^{-1}, \Sigma^{\frac{1}{2}}$ and $\Sigma^{-\frac{1}{2}}$.

$$(2\pi)^{-\frac{n}{2}}(\det\Sigma)^{-\frac{1}{2}}\int_{\mathbb{R}^{n}}e^{ix^{T}\xi}e^{-\frac{1}{2}(x-m)^{T}\Sigma^{-1}(x-m)}\;dx$$

$$=\qquad (2\pi)^{-\frac{n}{2}}(\det\Sigma)^{-\frac{1}{2}}\int_{\mathbb{R}^{n}}e^{ix^{T}\xi}e^{i(x-m)^{T}\xi}e^{-\frac{1}{2}(x-m)^{T}\Sigma^{-1}(x-m)}\;dx$$

$$=\qquad (2\pi)^{-\frac{n}{2}}(\det\Sigma)^{-\frac{1}{2}}e^{im^{T}\xi}\int_{\mathbb{R}^{n}}e^{i(x-m)^{T}\xi}e^{i(x-m)^{T}\xi}e^{-\frac{1}{2}(x-m)^{T}\Sigma^{-1}(x-m)}\;dx$$

$$y=\Sigma^{-\frac{1}{2}}x\qquad (2\pi)^{-\frac{n}{2}}e^{im^{T}\xi}\int_{\mathbb{R}^{n}}e^{i(\Sigma^{\frac{1}{2}}y)\xi}e^{-\frac{1}{2}|y|^{2}}\;dy$$

$$=\qquad (2\pi)^{-\frac{n}{2}}e^{im^{T}\xi}\int_{\mathbb{R}^{n}}e^{iy^{T}(\Sigma^{\frac{1}{2}}\xi)}e^{-\frac{1}{2}|y|^{2}}\;dy$$
Fourier transformation
$$=e^{im^{T}\xi}e^{-\frac{1}{2}|\Sigma^{\frac{1}{2}}\xi|^{2}}$$

$$=\qquad e^{im^{T}\xi}e^{-\frac{1}{2}\xi^{T}\Sigma\xi}$$

In terms of the uniqueness theorem of characteristic function (in [1], p.199, Satz 23.4), then we can deduce (2.4) is density function of X.

THEOREM 2.17. A linear combination of independent normally distributed random variable (or vector) is Gaussian.

Proof. We suppose X_1, \dots, X_m are independent random vectors on \mathbb{R}^n and $c_1, \dots, c_m \in$

 \mathbb{R} . Let us have a look at the characteristic function of it,

$$\begin{split} \mathbf{E}e^{i\xi^{T}\sum_{j=1}^{m}(c_{j}X_{j})} & \stackrel{independent}{=} & \prod_{j=1}^{m} \mathbf{E}e^{i\xi^{T}(c_{j}X_{j})} \\ & = & \prod_{j=1}^{m} \exp\left(i\xi^{T}\mathbf{E}[c_{j}X_{j}] - \frac{1}{2}\xi^{T}\mathrm{Var}[c_{j}X_{j}]\xi\right) \\ & = & \exp\left(i\xi^{T}\mathbf{E}[\sum_{j=1}^{m}c_{j}X_{j}] - \frac{1}{2}\xi^{T}\sum_{j=1}^{m}Var[c_{j}X_{j}]\xi\right) \\ & \stackrel{independent}{=} & \exp\left(i\xi^{T}\mathbf{E}[\sum_{j=1}^{m}c_{j}X_{j}] - \frac{1}{2}\xi^{T}\mathrm{Var}[\sum_{j=1}^{m}c_{j}X_{j}]\xi\right), \end{split}$$

which is a form of characteristic function of normal distribution. That means $\sum_{j=1}^{m} c_j X_j$ is Gaussian.

EXAMPLE 2.18 (Bivariate Normal Distribution). Cf.[16], p.241, Example 8.6. Suppose S_1, S_2 are independent random variables and have standard normal distributions. $\begin{pmatrix} S_1 \\ S_2 \end{pmatrix}$ has standard normal joint distribution since they are independent. We define

$$\begin{pmatrix} Y_1 \\ Y_2 \end{pmatrix} = \begin{pmatrix} \sigma_1, & 0 \\ \sigma_2 \rho, \sigma_2 (1 - \rho^2)^{\frac{1}{2}} \end{pmatrix} \cdot \begin{pmatrix} S_1 \\ S_2 \end{pmatrix} + \begin{pmatrix} \mu_1 \\ \mu_2 \end{pmatrix}, \tag{2.5}$$

where $\mu_1, \mu_2, \sigma_1, \sigma_2 \in \mathbb{R}, -1 \le \rho \le 1$. Again, Y_1, Y_2 are Gaussian and the joint distribution $\begin{pmatrix} Y_1 \\ Y_2 \end{pmatrix}$ is also Gaussian. Since S_1, S_2 are independent,

$$Var[Y_{1}] = Var[\sigma_{1}S_{1}]$$

$$= \sigma_{1}^{2},$$

$$Var[Y_{2}] = Var[\sigma_{2}\rho S_{1}] + Var[\sigma_{2}(1-\rho^{2})^{\frac{1}{2}}S_{2}]$$

$$= \sigma_{2}^{2}\rho^{2} + \sigma_{2}^{2}(1-\rho^{2})$$

$$= \sigma_{2}^{2},$$

$$Cov[Y_{1}, Y_{2}] = E[(Y_{1} - E[Y_{1}])(Y_{2} - E[Y_{2}])]$$

$$= E[Y_{1}Y_{2} - \mu_{1}Y_{2} - \mu_{2}Y_{1} + \mu_{1}\mu_{2}]$$

$$= E[(\sigma_{1}S_{1} + \mu_{1})(\sigma_{2}\rho S_{1} + \sigma_{2}(1-\rho^{2})^{\frac{1}{2}}S_{2} + \mu_{2})] - \mu_{1}\mu_{2}$$

$$= \sigma_{1}\sigma_{2}\underbrace{E[S_{1}^{2}]}_{=1}\rho + \mu_{1}\sigma_{2}\rho\underbrace{E[S_{1}]}_{=0} + \sigma_{1}\sigma_{2}(1-\rho^{2})^{\frac{1}{2}}\underbrace{E[S_{1}S_{2}]}_{=E[S_{1}]E[S_{2}]=0}$$

$$+ \mu_{1}\sigma_{2}(1-\rho^{2})^{\frac{1}{2}}\underbrace{E[S_{2}]}_{=0} + \sigma_{1}\underbrace{E[S_{1}]}_{=0}\mu_{2} + \mu_{1}\mu_{2} - \mu_{1}\mu_{2}$$

$$= \rho\sigma_{1}\sigma_{2},$$

that means the corrlation of Y_1, Y_2 is ρ . Because of the (2.4), the joint density function

$$f_{Y_1,Y_2}(y_1,y_2) = (2\pi)^{-1}(\det(\Sigma))^{-\frac{1}{2}}\exp((y_1-\mu_1)\Sigma^{-1}(y_2-\mu_2)),$$

where
$$\Sigma = \begin{pmatrix} \sigma_1^2, & 0 \\ \sigma_2^2 \rho^2, \sigma_2^2 (1 - \rho^2) \end{pmatrix}$$

Indeed,

$$\det(\Sigma) = (1 - \rho^2)\sigma_1^2 \sigma_2^2$$

and

$$\Sigma^{-1} = \frac{\begin{pmatrix} \sigma_2^2 (1 - \rho^2), & 0\\ -\sigma_2^2 \rho, & \sigma_1^2 \end{pmatrix}}{(1 - \rho^2)\sigma_1^2 \sigma_2^2}.$$

Namely,

$$f_{Y_1,Y_2}(y_1,y_2) = \frac{1}{2\pi(1-\rho^2)^{\frac{1}{2}}\sigma_1\sigma_2} \exp\left(-\frac{1}{2(1-\rho^2)}(z_1^2 - 2\rho z_1 z_2 + z_2^2)\right)$$
(2.6)

where $z_1 = \frac{y_1 - \mu_1}{\sigma_1}, z_2 = \frac{y_2 - \mu_2}{\sigma_2}$.

COROLLARY 2.19. Let Y_1, Y_2 be \mathbb{R} -valued random variables and $\begin{pmatrix} Y_1 \\ Y_2 \end{pmatrix}$ has a joint normal distribution, then the conditional expected value of Y_2 given Y_1

$$E[Y_2|Y_1 = y_1] = E[Y_2] + \rho(y_1 - E[Y_1]) \frac{\sigma_2}{\sigma_1},$$

and the conditional variance of Y_2 given Y_2

$$Var[Y_2|Y_1 = y_1] = \sigma_2^2(1 - \rho^2).$$

Where σ_1, σ_2 are standard deviations of Y_1, Y_2 and ρ is the correlation of Y_1, Y_2 .

Proof. Recall (2.6), we can specify the joint density function if σ_1, σ_2, ρ are known. As result of this, $\begin{pmatrix} Y_1 \\ Y_2 \end{pmatrix}$ has a form of (2.5). Suppose S_1, S_2 are independent standard normal distributed random variables. Now we have

$$S_1 \sim \frac{(Y_1 - \mathbf{E}[Y_1])}{\sigma_1}$$

 $Y_2 \sim \sigma_2 \rho S_1 + \sigma_2 (1 - \rho^2)^{\frac{1}{2}} S_2 + \mathbf{E}[Y_2],$

more precisely,

$$Y_2 \sim \sigma_2 \rho \frac{(Y_1 - \mathrm{E}[Y_1])}{\sigma_1} + \sigma_2 (1 - \rho^2)^{\frac{1}{2}} S_2 + \mathrm{E}[Y_2].$$

Take expectation of both sides,

$$E[Y_2|Y_1 = y_1] = \sigma_2 \rho \frac{(y_1 - E[Y_1])}{\sigma_1} + E[Y_2]$$

Now consider

$$Var[Y_2|Y_1 = y_1] = E[(Y_2 - \mu_{Y_2|Y_1})^2 | Y_1 = y_1]$$

$$= \int_{-\infty}^{\infty} (y_2 - \mu_{Y_2|Y_1})^2 f_{Y_2|Y_1}(y_2, y_1) dy_2$$

$$= \int_{-\infty}^{\infty} \left[y_2 - \mu_2 - \frac{\rho \sigma_2}{\sigma_1} (y_1 - \mu_1) \right]^2 f_{Y_2|Y_1}(y_2, y_1) dy_2,$$

After multiplying both sides by the density function of Y_1 and integrating it by y_1 , we have

$$\int_{-\infty}^{\infty} \operatorname{Var}[Y_{2}|Y_{1} = y_{1}] f_{Y_{1}}(y_{1}) dy_{1}$$

$$= \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \left[y_{2} - \mu_{2} - \frac{\rho \sigma_{2}}{\sigma_{1}} (y_{1} - \mu_{1}) \right]^{2} \underbrace{f_{Y_{2}|Y_{1}}(y_{2}, y_{1}) f_{Y_{1}}(y_{1})}_{f_{Y_{1}, Y_{2}}(y_{1}, y_{2})} dy_{2} dy_{1}$$

$$\iff \operatorname{Var}[Y_{2}|Y_{1} = y_{1}] \underbrace{\int_{-\infty}^{\infty} f_{Y_{1}}(y_{1}) dy_{1}}_{1}$$

$$= \operatorname{E}\left[(Y_{2} - \mu_{2}) - (\frac{\rho \sigma_{2}}{\sigma_{1}})(Y_{1} - \mu_{1}) \right]^{2}$$

multiplying right side out, we see

$$\operatorname{Var}[Y_{2}|Y_{1} = y_{1}] = \underbrace{\operatorname{E}[(Y_{2} - \mu_{2})^{2}]}_{\sigma_{2}^{2}} - 2 \frac{\rho \sigma_{2}}{\sigma_{1}} \underbrace{\operatorname{E}[(Y_{1} - \mu_{1})(Y_{2} - \mu_{2})]}_{\rho \sigma_{1} \sigma_{2}} + \frac{\rho^{2} \sigma_{2}^{2}}{\sigma_{1}^{2}} \underbrace{\operatorname{E}[(Y_{1} - \mu_{1})^{2}]}_{\sigma_{1}^{2}} \\
= \sigma_{2}^{2} - 2\rho^{2} \sigma^{2} + \rho^{2} \sigma_{2}^{2} \\
= \sigma_{2}^{2} - \rho^{2} \sigma_{2}^{2} \\
= \sigma_{2}^{2} (1 - \rho^{2}).$$

THEOREM 2.20. Let X be a Gaussian random variable, then

$$E[\exp(\beta X)] = \exp(\beta \mu + \frac{1}{2}\beta^2 \sigma^2). \tag{2.7}$$

Where μ and σ are E[X] and Var[X] respectively.

Proof.

$$\begin{split} & \operatorname{E}[\exp(\beta X)] \\ & = \ (2\pi\sigma^2)^{-\frac{1}{2}} \int_{\mathbb{R}} \exp(\beta x) \exp\left(-\frac{(x-\mu)^2}{2\sigma^2}\right) dx \\ & = \ (2\pi\sigma^2)^{-\frac{1}{2}} \int_{\mathbb{R}} \exp(\beta x) \exp\left(-\frac{(x^2-2x\mu+\mu^2)}{2\sigma^2}\right) dx \\ & = \ (2\pi\sigma^2)^{-\frac{1}{2}} \int_{\mathbb{R}} \exp\left(-\frac{x^2-2(\beta\sigma^2+\mu)x+\mu^2}{2\sigma^2}\right) dx \\ & = \ (2\pi\sigma^2)^{-\frac{1}{2}} \int_{\mathbb{R}} \exp\left(-\frac{x^2-2(\beta\sigma^2+\mu)x+(\beta\sigma^2+\mu)^2-(\beta\sigma^2+\mu)^2+\mu^2}{2\sigma^2}\right) dx \\ & = \ (2\pi\sigma^2)^{-\frac{1}{2}} \int_{\mathbb{R}} \exp\left(-\frac{(x-(\beta\sigma^2+\mu))^2+\mu^2-(\beta\sigma^2+\mu)^2}{2\sigma^2}\right) dx \\ & = \ \exp\left(\frac{(\beta\sigma^2+\mu)^2-\mu^2}{2\sigma^2}\right) \underbrace{(2\pi\sigma^2)^{-\frac{1}{2}} \int_{\mathbb{R}} \exp\left(-\frac{(x-(\beta\sigma^2+\mu))^2}{2\sigma^2}\right) dx}_{1} \\ & = \ \exp\left(\frac{\beta^2\sigma^4+2\mu\beta\sigma^2}{2\sigma^2}\right) \\ & = \ \exp(\mu\beta+\frac{1}{2}\beta^2\sigma^2) \end{split}$$

2.3 Brownian Motion

The Brownian motion was first introduced by Bachelier in 1900 in his PhD thesis. Now we give the common definition of it.

DEFINITION 2.21. Let $(B_t)_{t\geq 0}$ be a \mathbb{R}^n -valued stochastic process. (B_t) is called *Brownian motion* if it satisfies the following conditions:

- (i) $B_0 = 0$ a.s. .
- (ii) $(B_{t_1} B_{t_0}), \dots, (B_{t_n} B_{t_{n_1}})$ are independent for $0 = t_0 < t_1 < \dots < t_n$ and $n \in \mathbb{N}$.
- (iii) $B_t B_s \sim B_{t-s}$, for $0 \le s \le t < \infty$.
- (iv) $B_t B_s \sim \mathcal{N}(0, t s)^{\otimes n}$.
- (v) B_t is continuous in t a.s. .

A usual saying for (ii) and (iii) is the Brownian motion has independent, stationary increments. In (iv), \mathcal{N} represent a random variable which has a normal distribution. B_t is normally distributed due to (ii). It is clear that the increments of Brownian motion is stationary.

PROPOSITION 2.22. Let (B_t) be \mathbb{R} -valued Brownian motion. Then the covariance of B_m, B_n for $m, n \geq 0$ is $m \wedge n$.

Proof. Without loss of generality, we assume that $m \geq n$, then

$$E[B_m B_n] = E[(B_m - B_n)B_n] + E[B_n^2]$$
$$= E[B_m - B_n]E[B_n] + n$$
$$= n.$$

PROPOSITION 2.23. Let (B_t) be \mathbb{R} -valued Brownian motion. Then $B_{cm} \sim c^{\frac{1}{2}} B_m$.

Proof. Because B_m is normal distributed for any m > 0, we then get

$$E[e^{i\xi B_{cm}}] = e^{-\frac{1}{2}cm\xi^{2}}$$

$$= e^{-\frac{1}{2}m(c^{\frac{1}{2}}\xi)^{2}}$$

$$= E[e^{i\xi c^{\frac{1}{2}}B_{m}}].$$

Theorem 2.24. A \mathbb{R} -valued Brownian motion is a Gaussian process.

Proof. The following idea using the independence of increments to prove the claim come from [19]. We choose $0 = t_0 < t_1 < \cdots < t_n$, for $n \in \mathbb{N}$. Define $V = (B_{t_1}, \dots, B_{t_n})^T$, $K = (B_{t_1}, \dots, B_{t_n})^T$

$$(B_{t_1} - B_{t_0}, \dots, B_{t_n} - B_{t_{n-1}})^T \text{ and } A = \begin{pmatrix} 1 & 0 & \cdots & 0 \\ 1 & 1 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 1 & 1 & \cdots & 1 \end{pmatrix}. \text{ Let us look at the characteristic}$$

function of V,

In Proposition 2.3, $(t_j \wedge t_h)_{t,h=1,\dots,n}$ is the covariance matrix of V and therefore it is symmetric and positive definit. The mean vector of it is zero, then we have proved that the characteristic function is a form of some normal distributed random vector, i.e., V is Gaussian.

Schilling gave in his lecture [19] the relationship between a one-dimensional Brownian motion and a n-dimensional Brownian motion. In fact, $(B_t^{(l)})_{l=1,\dots,n}$ is Brownian motion if and only if $B_t^{(l)}$ is Brownian motion and all of the components are independent. Using this independence and the theorem of Fubini in the characteristic function for high dimensional Brownian motion we can say a n-dimensional Brownian motion is also a Gaussian process.

DEFINITION 2.25. Let $(X_t)_{t\in T}$ be a stochastic process. $(Y_t)_{t\in T}$ is defined on the same probability space as $(X_t)_{t\in T}$ and said to be *modification* of $(X_t)_{t\in T}$, if

$$\mathcal{P}[X_t = Y_t] = 1 \quad \forall \quad t \in T.$$

THEOREM 2.26 (Kolmogorov Chentsov). Let $(X_t)_{t\geq 0}$ be a stochastic process on \mathbb{R}^n such that

$$[|X_j - X_k|^{\alpha}] \le c|j - k|^{1+\beta} \quad \forall \quad j, k \ge 0 \quad \text{and} \quad j \ne k,$$

for $\alpha, \beta > 0, c < \infty$. Then $(X_t)_t$ has a modification $(Y_t)_t$ with continuous sample path such that

$$E\left[\left(\frac{|Y_j - Y_k|}{|j - k|^{\gamma}}\right)^{\alpha}\right] < \infty$$

for all $\gamma \in (0, \frac{\beta}{\alpha})$.

Proof. See [13], p.519.

Lemma 2.27. Let $(B_t)_{t\geq 0}$ be Brownian motion. Then

$$E[B_t^{2k}] = (2k - 1)!!t^{2k}$$

for $k \in \mathbb{N}$.

Proof. Cf.[19]. Taking expectation of B_t^{2k} , we get

$$\begin{split} \mathrm{E}[B_t^{2k}] &= \frac{1}{\sqrt{2\pi t}} \int_{-\infty}^{\infty} x^{2k} e^{-\frac{x^2}{2t}} \; dx \\ x = & \frac{2^k t^k}{\sqrt{\pi}} \int_{0}^{\infty} y^{k - \frac{1}{2}} e^{-y} \; dy \\ &= \frac{2^k t^k}{\sqrt{\pi}} \int_{0}^{\infty} y^{k + \frac{1}{2} - 1} e^{-y} \; dy \\ &= \frac{2^k t^k}{\sqrt{\pi}} \Gamma(k + \frac{1}{2}) \\ &= \frac{2^k t^k}{\sqrt{\pi}} \Gamma(\frac{1}{2}) \prod_{j=1}^k (j - \frac{1}{2}) \\ &= 2^k t^k \prod_{j=1}^k (\frac{2j-1}{2}) \\ &= (2k-1)!! \cdot t^k \end{split}$$

COROLLARY 2.28. Let $(B_t)_{t\geq 0}$ be Brownian motion. Then B_t is γ -Hölder continuous on a compact scale almost surely for all $\gamma < \frac{1}{2}$.

Proof. Because of Lemma 2.27, we have

$$E[(B_t - B_s)^{2k}] = E[B_{t-s}^{2k}]$$

= $(2k-1)!! \cdot |t-s|^k$.

In	${\rm terms}$	of	the	${\bf Theorem}$	of	${\bf Kolmogorov}$	Chenstov,	B_t	is	$\gamma\text{-H\"{o}lder}$	continuous	a.s.	for
γ ($\in (0, \frac{k}{2k})$).											

3 Stable Measures and Stable Integrals

In order to represent an integration form of fractional Brownian motion, we deal with the stable integral in this section. In fact, fractional Brownian motion is a Gaussian process with zero mean. To show Gaussian properties of it, we define it by a stable integral which can imaged as stochastic process of stable variables on time.

3.1 Stable Variables

DEFINITION 3.1. Let X be a random variable. X is said to have a stable distribution, if there exist $0 < \gamma \le 2, \delta \ge 0, -1 \ge \kappa \ge 1, \theta \in \mathbb{R}$ such that its characteristic function can be described as following

$$E[\exp i\xi X] = \begin{cases} \exp\{i\xi\theta - |\delta\xi|^{\gamma}(1 - i\kappa \cdot sgn(\xi)\tan\frac{\gamma\pi}{2})\}, & \text{if } \gamma \in (0,1) \cup (1,2], \\ \exp\{i\xi\theta - |\delta\xi|(1 + i\frac{2}{\pi}\kappa \cdot sgn(\xi)\ln|\xi|)\}, & \text{if } \gamma = 1. \end{cases}$$
(3.1)

Where

$$sgn(x) = \begin{cases} 1, & \text{if } x > 1 \\ 0, & \text{if } x = 0 \\ -1, & \text{if } x < 0. \end{cases}$$

Notice, we write $\Lambda(\gamma, \kappa, \theta, \delta)$ for one random variable whose characteristic function equals (3.1).

THEOREM 3.2. X is Gaussian if and only $X \sim \Lambda(\gamma, \kappa, \theta, \delta)$ with $\gamma = 2$.

Proof. On the one hand, if X is Gaussian, according to the charateric function of a Gaussian variable, γ must equal 2. On the other hand, if $\gamma = 2$, then $i\kappa \cdot sgn(\xi)\tan\frac{\gamma\pi}{2}$ vanishes since $\tan(\pi) = 0$. Therefore, X is Gaussian because $\mathrm{E}[\exp i\xi X] = \exp\{i\xi\theta - |\delta\xi|^2\}$.

Remark, if $\gamma = 2$, then κ is irrelevant in Definition. We specific $\kappa = 0$ without loss of generality. For instance, $B_t \sim \Lambda(2, 0, 0, \frac{\sqrt{2t}}{2})$ when $(B_t)_t$ is Brownian motion.

DEFINITION 3.3. A random variable X is said to be *symmetric* if X and -X have the same distribution.

PROPOSITION 3.4. Let X be have a stable distribution. X is *symmetric* if and only if $X \sim \Lambda(\gamma, 0, 0, \delta)$. I.e. its characteristic function has the form

$$E[\exp\{i\xi X\}] = \exp\{-|\delta\xi|^{\gamma}\} \tag{3.2}$$

Proof. The Definition of symmetricity implies

$$\exp\{i\xi\theta - |\delta\xi|^{\gamma}(1 - i\kappa \cdot sgn(\xi)\tan\frac{\gamma\pi}{2})\}$$

$$= \operatorname{E}[i\xi X]$$

$$= \operatorname{E}[i\xi(-X)]$$

$$= \operatorname{E}[i(-\xi)X]$$

$$= \exp\{i(-\xi)\theta - |\delta\xi|^{\gamma}(1 - i\kappa \cdot sgn(-\xi)\tan\frac{\gamma\pi}{2})\},$$

for $\xi \in \mathbb{R}$. This requires $\theta = \kappa = 0$.

COROLLARY 3.5. Let $(B_t)_t$ be Brownian motion, then B_t has a symmetric stable distribution.

Proof. It is clear due to the previous Proposition.

3.2 Stable Random Measures

In this subsection we suppose $(\Omega, \mathcal{A}, \mathcal{P})$, (D, \mathcal{D}, μ) are probability spaces, $\kappa(\cdot) : D \to [-1, 1]$ is a measurable function. For the next definition we need a notation

$$\mathscr{G} = \{ D \in \mathscr{D} : \mu[D] < \infty \}. \tag{3.3}$$

DEFINITION 3.6. Let ν be a set function such that

$$\nu: \mathscr{G} \to \mathcal{L}^0(\Omega).$$

 ν is said to be independently scattered, if $\nu[D_1], \ldots, \nu[D_n]$ are independent for any D_1, \ldots, D_n disjoint $\in \mathcal{D}$.

DEFINITION 3.7. Let ν be an independent cattered and σ -additive set function, ν is said to be *stable random measure* on (D, \mathcal{D}) with control measure μ , degree γ and skewness intensity $\kappa(\cdot)$ if

$$\nu[F] \sim \Lambda\left(\gamma, \frac{\int_{F} \kappa(x) \,\mu[dx]}{\mu[F]}, 0, (\mu[F])^{\frac{1}{\gamma}}\right) \tag{3.4}$$

for $F \in \mathscr{D}$.

Samorodnitsky and Tagqu showed the existence of stable measures, see [15], pp.119~120.

EXAMPLE 3.8. Suppose [0,T] is a index set, $0 = t_0 < t_1 < \cdots < t_k = T$ for $k \in \mathbb{N}$ and (B_t) is Brownian motion. We show the mapping $\nu : \mathcal{B}([0,T]) \to \mathcal{L}^0(\Omega)$, where $\nu[A_j](\omega) := B_{t_{j+1}}(\omega) - B_{t_j}(\omega)$ for $A_j = [t_j, t_{j+1})$.

Firstly, we show ν is independently scattered and σ -additive. We take $\{A_j\}$ such that $\bigcup_{j=1}^{\infty} A_j = [0,T]$. $\{\nu[A_k]\}_{k=1}^{\infty}$ has independent elements since $B_{t_1} - B_{t_0}, \ldots, B_{t_{j+1}} - B_{t_j}$ are independent.

Secondly,

$$\nu[\left(\bigcup_{j=1}^{\infty} A_j\right)] = B_T - B_1$$

$$= \sum_{j=1}^{\infty} (B_{t_{j+1}} - B_{t_j})$$

$$= \sum_{j=1}^{\infty} \nu[A_j].$$

Finally,

$$E[\exp(i\xi\nu[A_j])] = E[\exp(i\xi(B_{t_{j+1}} - B_{t_j})]$$

$$= \exp(-\frac{(t_{j+1} - t_j)\xi^2}{2})$$

Comparing with (3.1), we deduce the control measure must be $\frac{|\cdot|}{2}$. In fact, $\nu[A_j] \sim \Lambda(2,0,0,\frac{\sqrt{2|t_{j+1}-t_j|}}{2})$.

3.3 Stable Integrals

Samorodnitsky and Taqqu defined an integral with respect to stable measure as stochastic process in [15].

DEFINITION 3.9. The stable integral is given as follows

$$\int_{F} f(x) \,\nu(dx)(\omega) \tag{3.5}$$

with $f: F \to \mathbb{R}$ is a measurable function, given $\gamma \in (0,2], \mu: \mathcal{D} \to \mathcal{B}(\mathbb{R}), \kappa: \mathbb{R} \to \mathbb{R}$, such that

$$\begin{cases} \int_{F} |f(x)|^{\gamma} \mu(dx) < \infty, & \text{if } \gamma \in (0,1) \cup (1,2], \\ \int_{F} |\kappa(x)f(x) \ln |f(x)| |\mu(dx) < \infty, & \text{if } \gamma = 1, \end{cases}$$

$$(3.6)$$

where γ, μ, κ are, respectively, degree, control measure and skewness intensity of the stable measure ν .

Some properties of the stable function are given by Samorodnitsky and Taqqu.

Proposition 3.10. Let J(f) be a stable integral as form of (3.5). Then

$$J(f) \sim \Lambda(\gamma, \kappa, \theta, \delta)$$

with degree, control measure, skewness intensity, respectively,

$$\gamma \in (0,2],$$

$$\kappa = \frac{\int_{F} \kappa(x)|f(x)|^{\gamma} \cdot sgn(f(x))\mu(dx)}{\int_{F} |f(x)|^{\gamma}\mu(dx)},$$

$$\theta = \begin{cases} 0, & \text{if } \gamma \in (0,1) \cup (1,2], \\ -\frac{2}{\pi} \int_{F} \kappa(x)f(x) \ln |f(x)|\mu(dx), & \text{if } \gamma = 1, \end{cases}$$

$$\delta = \left(\int_{F} |f(x)^{\gamma}|\mu(dx)\right)^{\frac{1}{\gamma}},$$

of the stable measure ν .

Proof. See.[15], p.124, Proposition 3.4.1.

Proposition 3.11. The stable integral is linear, in fact,

$$J(c_1f_1 + c_2f_2) \stackrel{a.s.}{=} c_1J(f_1) + c_2J(f_2)$$
(3.7)

for any f_1, f_2 integrable with respect to some stable measure and real numbers c_1, c_2 .

Proof. See [15], p.117, Property
$$3.2.3$$
.

4 Fractional Brownian Motion

The fractional Brownian motion (fBm) was defined by Kolmogorov primitively. After that Mandelbrot and Van Ness has presented the work in detail. This section is concerned with the definition and some properties of it.

4.1 Definition of Fractional Brownian Motion

Mandelbrot and Van Ness [17] gave an integration representation of fBm.

DEFINITION 4.1. Let $(U_H(t))_{t\in\mathbb{R}}$ be a \mathbb{R} -valued stochastic process and H be a real number such that 0 < H < 1. $(U_H(t))$ is said to be fractional Brownian motion if

$$U_H(t) - U_H(s) = \frac{1}{\Gamma(H + \frac{1}{2})} \left(\int_{\mathbb{R}} \mathbb{1}_{\{t > u\}} \cdot (t - u)^{H - \frac{1}{2}} - \mathbb{1}_{\{s > u\}} (-u)^{H - \frac{1}{2}} dB_u \right) (4.1)$$

for $t \geq s, t, s \in \mathbb{R}$. Where (B_u) is defined as two-sides Brownian motion and the integral is defined in sense of stable integral as in previous section. H is called $Hurst\ exponent$ or $Hurst\ index$ of fBm.

As usual, we set $U_H(0) = 0$, then equation (4.1) is equivalent to

$$U_H(t) = \frac{1}{\Gamma(H + \frac{1}{2})} \left(\int_{\mathbb{R}} \mathbb{1}_{\{t > u\}} \cdot (t - u)^{H - \frac{1}{2}} - \mathbb{1}_{\{u < 0\}} (-u)^{H - \frac{1}{2}} dB_u \right). \tag{4.2}$$

Lemma 4.2. The equation (4.2) is well-defined, $U_H(t)$ has stable distribution and

$$U_H(t) \sim \Lambda(2, 0, 0, \frac{1}{\Gamma(H + \frac{1}{2})} \left(\int_{\mathbb{R}} |f(u)^2| \frac{du}{2} \right)^{\frac{1}{2}},$$

Where f(x) is the integrand of integral in (4.2).

Proof. Firstly, B_t is Gaussian and symmetric stable measure with zero mean and $\frac{|\cdot|}{2}$ is the control measure of it shown in Example 3.8.

Secondly, by $H = \frac{1}{2}$, $\int_{\mathbb{R}} f^2 \frac{du}{2} = \frac{1}{2} \int_0^{|t|} du = \frac{1}{2}|t| < \infty$. By $H \neq \frac{1}{2}$, we deal it with Taylor expansion. As u goes to $-\infty$, $f(u) = -(H - \frac{1}{2})(t - u)^{H - \frac{3}{2}} - (H - \frac{1}{2})(-u)^{H - \frac{3}{2}} + o(1)$. Where o(1) tends to zero when u is around $-\infty$. Consider $H - \frac{3}{2} < 0$, then f(u) is square integrable around $-\infty$. As u goes to t, $f(u) \propto \mathbbm{1}_{\{t>u\}}(t-u)^{H - \frac{1}{2}}$. Hence, f(u) is also square integrable around u = t. It is clear f(u) = 0 when u is around ∞ . Then it satisfies the condition $\int_{-\infty}^{\infty} f^2(u) \frac{du}{2} < \infty$.

Finally, in terms of Proposition 3.10, we get the claim.

It is worth mentioning that, if we take $H = \frac{1}{2}$ and choose a restriction of the integrand on \mathbb{R}_+ , i.e., $U_{\frac{1}{2}}(t) = \frac{1}{\Gamma(H+\frac{1}{2})}B_t$ is a Brownian motion.

LEMMA 4.3. Let $(U_H(t))_t$ be a fBm. Then $U_H(t) \sim \mathcal{N}(0, \frac{1}{\Gamma(H+\frac{1}{2})^2}(\int_{\mathbb{R}} |f(u)^2| \ du))$.

Proof. In terms of Lemma (4.2), $E[i\xi U_H(t)] = \exp\{-\xi^2 \frac{1}{2\Gamma(H+\frac{1}{2})^2} (\int_{\mathbb{R}} |f(u)^2| \ du)\}$. The rest is clear thanks to the form of characteristic function of a Gaussian random variable. \Box

Notice that we can also define $U_H(t)$ by (4.2) with Itó integral. It has the same expected value and variance as defined by stable integral. Since $U_H(t)$ is Gaussian, both of two versions have the same distribution. Following properties remains true also by Itó integral version.

Corollary 4.4.
$$U_H(t) - U_H(s) \sim \mathcal{N}(0, \frac{1}{\Gamma(H + \frac{1}{2})^2} (\int_{\mathbb{R}} |f_t(u) - f_s(u)|^2 du)$$

Proof. This Corollary follows Proposition 3.7 and Lemma 4.3.

THEOREM 4.5. Let $(U_H(t))_t$ be a fBm. Then $U_H(t)$ has an expected value 0 and variance $\frac{1}{(\Gamma(H+\frac{1}{2}))^2}t^{2H}$ $EU_H^2(1)$ for any $t \in \mathbb{R}$.

Proof. It is clear that U_H is Gaussian with zero mean due to Lemma 4.2. We suppose that $t \geq s \geq 0, c(H) = \frac{1}{(\Gamma(H+\frac{1}{2}))^2}$.

$$\begin{aligned}
& \operatorname{E}[(U_{H}(t) - U_{H}(s))^{2}] \\
& = c(H)\operatorname{E}[\int_{\mathbb{R}} \left(\mathbb{1}_{\{t>u\}} \cdot (t-u)^{H-\frac{1}{2}} - \mathbb{1}_{\{s>u\}} \cdot (s-u)^{H-\frac{1}{2}} \right)^{2} du] \\
& = c(H)\operatorname{E}[\int_{\mathbb{R}} \left(\mathbb{1}_{\{t-s>u\}} \cdot (t-s-u)^{H-\frac{1}{2}} - \mathbb{1}_{\{0>u\}} \cdot (-u)^{H-\frac{1}{2}} \right)^{2} du] \\
& \stackrel{m=t-s}{=} c(H)\operatorname{E}[\int_{\mathbb{R}} \left(\mathbb{1}_{\{m>u\}} \cdot (m-u)^{H-\frac{1}{2}} - \mathbb{1}_{\{0>u\}} \cdot (-u)^{H-\frac{1}{2}} \right)^{2} du] \\
& \stackrel{u=ml}{=} c(H)\operatorname{E}[\int_{\mathbb{R}} \left(\mathbb{1}_{\{m>ml\}} \cdot (m-ml)^{H-\frac{1}{2}} - \mathbb{1}_{\{0>ml\}} \cdot (-ml)^{H-\frac{1}{2}} \right)^{2} m \cdot dl] \\
& = c(H)\operatorname{E}[\int_{\mathbb{R}} \left(\mathbb{1}_{\{1>l\}} \cdot (1-l)^{H-\frac{1}{2}} - \mathbb{1}_{\{0>l\}} \cdot (-l)^{H-\frac{1}{2}} \right)^{2} \cdot m^{2H-1} \cdot m \, dl] \\
& = c(H)m^{2H}\operatorname{E}[U_{H}(1)^{2}] \\
& = c(H)(t-s)^{2H}\operatorname{E}[U_{H}(1)^{2}]
\end{aligned} \tag{4.3}$$

Using the same calculation, we get

$$E[(U_H(t)^2] = c(H)t^{2H}E[U_H(1)^2].$$
(4.4)

(4.4) is variance of
$$U_H(t)$$
 due to $E[(U_H(t))] = 0$.

In order to normalize the variance, a definition of standard fBm is given.

DEFINITION 4.6. A stochastic process $(U_H(t))_t$ is said to be a standard fractional Brownian motion (sfBm) if

$$U_H(t) = \hat{c}(H) \int_{\mathbb{R}} \mathbb{1}_{\{t \ge u\}} \cdot (t - u)^{H - \frac{1}{2}} - \mathbb{1}_{\{u \le 0\}} (-u)^{H - \frac{1}{2}} dB_u.$$
 (4.5)

Where $\hat{c}(H) = \frac{1}{(\Gamma(H + \frac{1}{2})^2) E[U_H(1)^2]}$.

We consider from now on sfBm instead of fBm.

THEOREM 4.7. Let $(U_H(t))_t$ be a fBm. The covariance of $U_H(t)$ and $U_H(s)$ is $\frac{1}{2}(t^{2H} + s^{2H} - |t - s|^{2H})$ for $t, s \in \mathbb{R}$.

Proof. Cf.[17], Theorem 5.3.

$$Cov[U_{H}(t), U_{H}(s)] = E[U_{H}(t)U_{H}(s)]$$

$$= \frac{1}{2} \left(E[U_{H}(t)^{2}] + E[U_{H}(s)^{2}] - E[(U_{H}(t) - U_{H}(s))^{2}] \right)$$

$$\stackrel{(4.4)}{=} \frac{1}{2} (t^{2H} + s^{2H} - |t - s|^{2H})$$

$$(4.6)$$

THEOREM 4.8. $(U_H(t))_t$ is Gaussian process.

Proof. We just need to prove that for any finite linear combination of $(U_H(t))_t$ is Gaussian. We take $t_1, \ldots, t_k \in \mathbb{R}, c_1, \ldots, c_k \in \mathbb{R}$ and the stable integral J(f) is a linear functional with $\gamma = 2, \kappa = 0, \theta = 0, \delta = (\frac{1}{2} \int_{-\infty}^{\infty} f^2(u) \ du)^{\frac{1}{2}}$ due to Corollary 3.7. Suppose f_1, \ldots, f_k are integrands of stable integration of $U_H(t_1), \ldots, U_H(t_k)$ respectively.

Consider now, according to the Minkowski inequality,

$$\int_{-\infty}^{\infty} (\sum_{j=1}^{k} c_j f_j)^2 du \leq \sum_{j=1}^{k} \underbrace{\int_{-\infty}^{\infty} (c_j f_j)^2 du}_{<\infty}$$

$$< \infty.$$

Moreover,

$$\sum_{j=1}^{k} c_{j} U_{H}(t_{j}) = \sum_{j=1}^{k} c_{j} J(f_{j})$$

$$= J(\sum_{j=1}^{k} c_{j} f_{j})$$

$$\sim \Lambda(2, 0, 0, (\frac{1}{2} \int_{-\infty}^{\infty} (\sum_{j=1}^{k} c_{j} f_{j})^{2} du)^{\frac{1}{2}})$$

is Gaussian and the rest follows from Corollary 2.15.

COROLLARY 4.9. Let $(U_H(t))_t$ be a fBm, then $(U_H(t))_t$ has stationary and H-self similar increments.

Proof. Assume that $s \geq u$. Because the joint distribution of $(U_H(s), U_H(u))^T$ is Gaussian, $(1,-1)\cdot (U_H(s), U_H(u))^T$ is also Gaussian. From (4.3), $U_H(t_k+\tau) - U_H(s_k+\tau) \sim U_H(t_k) - U_H(s_k) \sim \mathcal{N}(0, (t_k-s_k)^{2H})$ for $k \in \{1 \dots d\}$. Corresponding to (2.3),

$$E[i\sum_{k=1}^{d} \xi_k (U_H(t_k + \tau) - U_H(s_k + \tau))] = E[i\sum_{k=1}^{d} \xi_k (U_H(t_k) - U_H(s_k))]$$

due to $(U_H(t_{k+\tau} - s_{k+\tau}))_{k=1}^d$ and $(U_H(t_k - s_k))_{k=1}^d$ have the same expected vector and covariance matrix in their characteristic function, i.e., $(U_H(t))$ has stationary increments. In order to show fBm has H-self similar increments, we have to prove $(U_H(zt_1), U_H(zt_2), \dots, U_H(zt_n)) \sim (z^H U_H(t_1), z^H U_H(t_2), \dots, z^H U_H(t_n))$ for any z > 0. Obviously, the former and the latter of the term is Gaussian and $\text{Var}[U_H(zt_i), U_H(zt_j)] = \text{Var}[z^H U_H(t_i), z^H U_H(t_j)] = \frac{1}{2} z^{2H} (t_i^{2H} + t_j^{2H} - |t_i - t_j|^{2H})$. Thus they have the same expected vector and covariance matrix in their characteristic function. In the same way as mentioned above, we get our claim.

4.2 Regularity

THEOREM 4.10 (Kolmogorov Chentsov). FBM has almost surely continuous sample path.

Proof. Cf.[17], Proposition 4.1. Let $(U_H(t))_t$ be a FBM with Hurst index H. Fix α such that $1 < \alpha H$. Let us have a look at the expectation of $(U_H(t) - U_H(s))^{\alpha}$ with respect to the calculation in (4.3)

$$E[(U_{H}(t) - U_{H}(s))^{\alpha}] = |t - s|^{\alpha H} \cdot \underbrace{E\left(\int_{\mathbb{R}} \mathbb{1}_{\{1 > u\}} \cdot (1 - u)^{H - \frac{1}{2}} - \mathbb{1}_{\{u < 0\}} (-u)^{H - \frac{1}{2}} dB_{u}\right)^{\alpha}}_{c(\alpha, H)}$$

$$= c(\alpha, H) \cdot |t - s|^{\alpha H}. \tag{4.7}$$

We choose $\beta = \alpha H - 1$ and $\gamma \in (0, H - \frac{1}{\alpha})$ then the rest follows from Theorem 2.26 . \Box

Remark, $(U_H(t))_t$ is, in fact, γ -Hölder continuous with $\gamma < H$ almost surely.

THEOREM 4.11. The sample path of FBM is almost surely not differentiable.

Proof. Cf. [17] Proposition 4.2 . Fix $\omega \in \Omega$, we assume $c > 0, t_j \to s$.

$$\mathcal{P}[\limsup_{t \to s} \left| \frac{U_H(t) - U_H(s)}{t - s} \right| > c]$$

$$= \mathcal{P}[\limsup_{j \to \infty} \sup_{t_j \neq s} \left| \frac{U_H(t_j) - U_H(s)}{t_j - s} \right| > c]$$
(4.8)

Since continuity of measures from above, then

$$(4.8) = \lim_{j \to \infty} \mathcal{P}[\sup_{t_j \neq s} |\frac{U_H(t_j) - U_H(s)}{t_j - s}| > c]$$

$$\geq \lim_{j \to \infty} \mathcal{P}[|\frac{U_H(t_j) - U_H(s)}{t_j - s}| > c]$$

$$= \lim_{j \to \infty} \mathcal{P}[\frac{|(t_j - s)^H U_H(1)}{t_j - s}| > c]$$

$$= \lim_{j \to \infty} \mathcal{P}[|(t_j - s)^{H-1} U_H(1)| > c]$$

$$= \lim_{j \to \infty} \mathcal{P}[|U_H(1)| > \underbrace{|t_j - s|^{1-H}}_{j \to \infty} c]$$

$$\xrightarrow{j \to \infty} 1$$

THEOREM 4.12. Let $(U_H(k))_k$ be a fBm. The conditional expectation of $U_H(s)$ given $U_H(t) = x$ is

$$\frac{|\frac{s}{t}|^{2H} + 1 - |\frac{s}{t} - 1|^{2H}}{2} \cdot x$$

for all $s \ge t$ and $t \ne 0$.

Proof. Cf. [17] Theorem 5.3. Taking conditional expectation of $U_H(s)$ given $U_H(t)$,

4.3 Fractional Brownian Noise

DEFINITION 4.13. Let $(U_H(t))_{t\in\mathbb{R}}$ be a fBm. The fractional Brownian noise is a sequence $(S_k)_{k\in\mathbb{R}}$ defined as follows

$$S_H(k) = U_H(k+1) - U_H(k)$$

for $k \in \mathbb{R}$.

Proposition 4.14. Fractional Brownian noise is stationary and its autocovariance is

$$\varsigma_{S_H}(\tau) = \frac{1}{2}(|\tau+1|^{2H} - 2|\tau|^{2H} + |\tau-1|^{2H})$$
(4.9)

for $\tau \in \mathbb{R}$.

Proof. Cf. [18], p.333, Proposition 7.2.9.

The first part of the claim is clear due to fBm has stationary increments. In terms of definition of fractional Brownian noise, for a $k \in \mathbb{R}$, we have

$$\varsigma_{S_{H}}(\tau)
= E[S_{H}(k+\tau)S_{H}(k)]
= E[(U_{H}(\tau+k+1) - U_{H}(\tau+k))(U_{H}(k+1) - U_{H}(k))]
= E[U_{H}(\tau+k+1)U_{H}(k+1)] - E[U_{H}(\tau+k+1)U_{H}(k)]
- E[U_{H}(\tau+k)U_{H}(k+1)] + E[U_{H}(\tau+k)U_{H}(k)]
\stackrel{(4.6)}{=} \frac{1}{2}((\tau+k+1)^{2H} + (k+1)^{2H} - |\tau|^{2H} - (\tau+k+1)^{2H} - k^{2H} + |\tau+1|^{2H}
- (\tau+k)^{2H} - (k+1)^{2H} + |\tau-1|^{2H} + (\tau+k)^{2H} + k^{2H} - |\tau|^{2H})
= \frac{1}{2}(|1+\tau|^{2H} - 2|\tau|^{2H} + |1-\tau|^{2H})$$

for $\tau \in \mathbb{R}$.

DEFINITION 4.15. A stationary stochastic process $(X_t)_t$ is said to have *long memory* if its autocovariance $\varsigma_X(\tau)$ tends to 0 so slowly such that $\sum_{\tau=0}^{\infty} \varsigma_X(\tau)$ diverges.

Lemma 4.16 (Cauchy Condensation test). Let (a_n) be a \mathbb{R} -valued positive non-increasing sequence. Then $\sum_{n=1}^{\infty} a_n$ converges if and only if $\sum_{n=0}^{\infty} 2^n a_{2^n}$ converges.

Proof. See [14], p.391, Theorem 13.13.

LEMMA 4.17 (Limit comparison test). Let $\sum_{k=0}^{\infty} a_k$ and $\sum_{k=0}^{\infty} b_k$ be two series with $a_k \geq 0, b_k \geq 0$ for all k. If $0 < \lim_{k \to \infty} \frac{a_k}{b_k} < \infty$, then either both series converge or both series diverge.

Proof. Suppose $\lim_{k\to\infty}\frac{a_k}{b_k}=c$ with $c<\infty$. Then there exists a N such that if $k>N, |\frac{a_k}{b_k}-c|<\frac{c}{2}$. In other words,

$$\frac{c}{2}|b_k| < |a_k| < \frac{3c}{2}|b_k|$$

Hence, if $\sum_{k=0}^{\infty} a_k$ converges, then $\sum_{k=0}^{\infty} b_k$ converges due to the first '<'. Also, if $\sum_{k=0}^{\infty} b_k$ converges then $\sum_{k=0}^{\infty} a_k$ converges due to the last '<'. One can easy to verify the claim of the divergence in the same way.

Theorem 4.18. The fractional Brownian noise with $H \in (\frac{1}{2}, 1)$ has long memory.

Proof. Cf. [18], p.335, Proposition 7.2.10.

Without loss of generality, we suppose $\tau \in \mathbb{N}$ because $\varsigma_{S_H}(0) = 1$.

$$\begin{split} &\varsigma_{S_H}(\tau) \\ &= \frac{1}{2}\tau^{2H-2}\{\tau^2[(1+\frac{1}{\tau})^{2H}-2+(1-\frac{1}{\tau})^{2H}]\} \\ &= \frac{1}{2}\tau^{2H-2}\{\frac{(1+\frac{1}{\tau})^{2H}-1}{\frac{1}{\tau^2}}-\frac{1-(1-\frac{1}{\tau})^{2H}}{\frac{1}{\tau^2}}\} \end{split}$$

We deal with the former of the content in $\{\}$ with L'Hôpital's rule as τ tends to infinity.

$$\frac{(1+\frac{1}{\tau})^{2H}-1}{\frac{1}{\tau^2}}$$

$$= \frac{2H(1+\frac{1}{\tau})^{2H-1}(-\frac{1}{\tau^2})}{-\frac{2}{\tau^3}} + o(1)$$

$$= \frac{H(1+\frac{1}{\tau})^{2H-1}}{\frac{1}{\tau}} + o(1)$$

We calculate the Latter of the content in {} in a similar way. Then

$$\begin{split} &\varsigma_{S_H}(\tau) \\ &= \quad \tau^{2H-2}\frac{1}{2}\{\frac{H(1+\frac{1}{\tau})^{2H-1}}{\frac{1}{\tau}} - \frac{H(1-\frac{1}{\tau})^{2H-1}}{\frac{1}{\tau}}\} + \mathrm{o}(1) \\ &= \quad \tau^{2H-2}\frac{1}{2}\{\frac{H(1+\frac{1}{\tau})^{2H-1} - H}{\frac{1}{\tau}} - \frac{H(1-\frac{1}{\tau})^{2H-1} - H}{\frac{1}{\tau}}\} + \mathrm{o}(1) \\ ^{\mathrm{L'Hopital}} &= \quad \frac{1}{2}\tau^{2H-2}\{\frac{H(2H-1)(1+\frac{1}{\tau})^{2H-2}(-\frac{1}{\tau^2})}{-\frac{1}{\tau^2}} - \frac{H(2H-1)(1-\frac{1}{\tau})^{2H-2}\frac{1}{\tau^2}}{-\frac{1}{\tau^2}}\} + \mathrm{o}(1) \\ &= \quad \frac{1}{2}\tau^{2H-2}\{H(2H-1)(1+\frac{1}{\tau})^{2H-2} + H(2H-1)(1-\frac{1}{\tau})^{2H-2}\} + \mathrm{o}(1) \\ &= \quad \frac{1}{2}\tau^{2H-2}2H(2H-1) + \mathrm{o}(1) \\ &= \quad H(2H-1)\tau^{2H-2} + \mathrm{o}(1) \end{split}$$

When $H \in (0, \frac{1}{2}), \sum_{\tau=1}^{\infty} \tau^{2H-2}$ converges. We use the Cauchy condesation test to verify it

$$\sum_{\tau=0}^{\infty} 2^{\tau} (2^{\tau})^{2H-2}$$

$$= \sum_{\tau=0}^{\infty} 2^{\tau(2H-1)}$$

This is geometric series if 2H-1<0, namely, $H\in(0,\frac{1}{2})$.

Otherwise, if $\frac{1}{2} < H < 1$, namely, -1 < 2H - 2 < 0, $\sum_{\tau=1}^{\infty} \tau^{2H-2}$ diverges because it is greater than the harmolic series.

H(2H-1) > 0 when $H \in (\frac{1}{2}, 1)$. Since the limit comparison test, $\sum_{\tau=1}^{\infty} \varsigma(\tau)$ diverges because $\sum_{\tau=1}^{\infty} \tau^{2H-2}$ diverges. It is clear that $\lim_{\tau \to \infty} \varsigma_{S_H}(\tau) = 0$, then claim is proved.

COROLLARY 4.19. Let S_H be fractional Brownian noise and $\varsigma_{S_H}(\cdot)$ be its autocovariance. Then $\sum_{\tau=0}^{\infty} \varsigma_{S_H}^2(\tau) < \infty$ if and only if $H < \frac{3}{4}$.

Proof. Cf. [18], p.72, Lemma 6.3. As by Theorem 4.18, we have $\varsigma_{S_H}^2(\tau) = H^2(2H-1)^2\tau^{4H-4} + o(1)$. The sum of it over the range of τ is finite if and only if, according to the same reason as in Theorem 4.18, 4H < 3. That means $H < \frac{3}{4}$.

4.4 fBm is not Semimartingale for $H \neq \frac{1}{2}$

Let us have a look at our integration representation for fBm. In the case of fBm with Hurst index $\frac{1}{2}$, it must be an ordinary Brownian motion. Otherwise, we'll show fBm is not a seimimartingale.

DEFINITION 4.20. The *Hermite polynomials* forms as following

$$H_n(u) = (-1)^n e^{\frac{u^2}{2}} (\frac{\partial^n}{\partial u^n} e^{-\frac{u^2}{2}}),$$
 (4.10)

for $u \in \mathbb{R}, n \in \mathbb{N}_0$.

PROPOSITION 4.21. Let $(H_n)_{n\in\mathbb{N}_0}$ be a family of Hermite polynomials, $f,g:\mathbb{R}\to\mathbb{R}$ continuously differentiable. It has then following properties

- (i) $H_{n+2}(u) = u \cdot H_{n+1}(u) (n+1)H_n(u)$ and $H_{n+1}(u) = (n+1)H'_n(u)$ for all $n \in \mathbb{N}_0, u \in \mathbb{R}$.
- (ii) Let W, V be standard Gaussian such that (W, V) have a disjoint Gaussian distribution. Then

$$\int_{\Omega} H_j(W) \cdot H_k(V) \mathcal{P} = \begin{cases} j! (\mathbf{E}[WV])^j & \text{if } j = k, \\ 0 & \text{otherwise.} \end{cases}$$

(iii) Let W be standard Gaussian distributed, then

$$\frac{1}{j!} \int_{\Omega} H_j(W) H_k(W) \mathcal{P} = \begin{cases} 1 & \text{if } j = k, \\ 0 & \text{otherwise.} \end{cases}$$

Remark that, (iii) means the fact, $\{\frac{1}{\sqrt{j!}}\cdot H_j(x)\}_{j=0}^{\infty}$ is an orthonormal basis in $\mathcal{L}^2(\mathbb{R}, \mathcal{B}(\mathbb{R}), e^{-\frac{x^2}{2}} dx)$.

Proof. See [18] p.3, Propostion 1.3.

LEMMA 4.22. Let $(U_H(t))_t$ be a fBm, W standard Gaussian variable, and $f: \mathbb{R} \to \mathbb{R}$, Borel-measurable function such that $\mathrm{E}[f^2(W)] < \infty$. Then,

$$\frac{1}{n} \sum_{j=1}^{n} f(U_H(j) - U_H(j-1)) \stackrel{\text{in } \mathcal{L}^2(\mathcal{P})}{\to} E[f(W)],$$

as n tends to ∞ . In particular,

$$\sum_{j=1}^{n} |U_{H}(\frac{j}{n}) - U_{H}(\frac{j-1}{n})|^{\beta} \stackrel{\text{in } \mathcal{L}^{2}(\mathcal{P})}{\longrightarrow} \begin{cases} 0 & \text{if } \beta > \frac{1}{H} \\ \mathbb{E}[|W|^{\beta}] & \text{if } \beta = \frac{1}{H} \\ \infty & \text{if } \beta < \frac{1}{H} \end{cases}$$

$$(4.11)$$

as n tends to ∞ .

Proof. C.f. [18], p.17, Theorem 2.1.

Firstly, because $\mathrm{E}[f^2(W)] < \infty$, one has $f \in \mathcal{L}^2(\mathbb{R}, \mathscr{B}(\mathbb{R}), e^{-\frac{x^2}{2}} dx)$. In terms of Proposition 4.21(iii), taking expectation

$$E[f(x)] = E[\sum_{j=0}^{\infty} \frac{a_j H_j(x)}{\sqrt{j!}}],$$

for $x \in \mathbb{R}$. Notice $H_0(u) = 1$ for $u \in \mathbb{R}$ due to (4.10). Setting x = W, equalling coefficients leads to $a_0 = \mathrm{E}[f(W)]$. Moreover,

$$E\left[\left\{\frac{1}{n}\sum_{j=1}^{n}f(U_{H}(j)-U_{H}(j-1))-E[f(W)]\right\}^{2}\right]$$

$$= E\left[\left\{\frac{1}{n}\sum_{j=1}^{n}(f(U_{H}(j)-U_{H}(j-1))-E[f(W)])\right\}^{2}\right]$$

$$= E\left[\left\{\frac{1}{n}\sum_{j=1}^{n}(\sum_{k=0}^{\infty}\frac{a_{k}}{\sqrt{k!}}H_{k}(U_{H}(j)-U_{H}(j-1)))-E[f(W)]\right\}^{2}\right]$$

$$= E\left[\left\{\frac{1}{n}\sum_{j=1}^{n}(\sum_{k=0}^{\infty}\frac{a_{k}}{\sqrt{k!}}H_{k}(U_{H}(j)-U_{H}(j-1)))\right\}^{2}\right]$$

Consider now

$$\mathrm{E}[f^2(W)] < \infty,$$

which requires $\sum_{k=1}^{\infty} (a_k)^2 < \infty$. Then

$$E[\left\{\frac{1}{n}\sum_{j=1}^{n}f(U_{H}(j)-U_{H}(j-1))-E[f(W)]\right\}^{2}]$$

$$=\frac{1}{n^{2}}E[\sum_{k=1}^{\infty}\frac{a_{k}^{2}}{k!}(\sum_{j=1}^{n}H_{k}(U_{H}(j)-U_{H}(j-1)))^{2}]$$

$$=\frac{1}{n^{2}}\sum_{k=1}^{\infty}\frac{a_{k}^{2}}{k!}\sum_{j=1,m=1}^{n}E[H_{k}(U_{H}(j)-U_{H}(j-1))H_{k}(U_{H}(m)-U_{H}(m-1)]$$
Proposition 4.21(ii)
$$=\frac{1}{n^{2}}\sum_{k=1}^{\infty}a_{k}^{2}\sum_{j=1,m=1}^{n}(E[(U_{H}(j)-U_{H}(j-1))(U_{H}(m)-U_{H}(m-1))])^{k}$$

$$=\frac{1}{n^{2}}\sum_{k=1}^{\infty}a_{k}^{2}\sum_{j=1,m=1}^{n}(E[S_{H}(j-1)S_{H}(m-1)])^{k}$$

$$=\frac{1}{n^{2}}\sum_{k=1}^{\infty}a_{k}^{2}\sum_{j=1,m=1}^{n}(\varsigma_{S_{H}}(j-m))^{k}$$

Notice that,

$$\begin{aligned} |\varsigma_{S_H}(k)| &= |\varsigma_{S_H}(|k|)| \\ &= & \mathrm{E}[(U_H(1) - U_H(0))(U_H(|x|+1) - U_H(|x|))] \\ &\leq & \underbrace{\sqrt{\mathrm{E}[U_H(1)^2]} \cdot \sqrt{\mathrm{E}[U_H(|k|+1) - U_H(|k|)]^2}}_{=1} \\ &= & 1 \end{aligned}$$

Consequently, $(\varsigma_{S_H}(j-m))^k \leq |\varsigma_{S_H}(j-m)|$. In fact,

$$E[\left\{\frac{1}{n}\sum_{j=1}^{n}f(U_{H}(j)-U_{H}(j-1))-E[f(W)]\right\}^{2}]$$

$$\leq \frac{1}{n^{2}}\sum_{\substack{k=1\\ =:\alpha<\infty}}^{\infty}a_{k}^{2}\sum_{j=1,m=1}^{n}|\varsigma_{S_{H}}(j-m)|$$

$$= \frac{\alpha}{n^{2}}\sum_{j=1,m=1}^{n}|\varsigma_{S_{H}}(j-m)|$$

$$= \frac{\alpha}{n^{2}}2\cdot\sum_{j=1}^{n}\sum_{m< j}|\varsigma_{S_{H}}(j-m)|$$

$$\leq \frac{\alpha}{n^{2}}2n\sum_{k=1}^{n-1}|\varsigma_{S_{H}}(k)|$$

$$= \frac{2\alpha}{n}\sum_{k=1}^{n-1}|\varsigma_{S_{H}}(k)|,$$

As in the proof in Theorem 4.18,

$$\begin{split} &\sum_{k=1}^{n-1} |\varsigma_{S_H}(k)| \\ &\propto & H(2H-1) \sum_{k=1}^{n-1} k^{2H-2} \\ &\propto & H(2H-1) n \cdot n^{2H-2} \\ &\propto & n^{2H-1}, \end{split}$$

as n goes to infinity. Then

$$\frac{2\alpha}{n} \sum_{k=1}^{n-1} |\varsigma_{S_H}(k)| \propto n^{2H-2}$$

as n goes to infinity. This leads to, for 0 < H < 1, as $n \to \infty$, $\mathrm{E}[\{\frac{1}{n}\sum_{j=1}^n f(U_H(j) - U_H(j-1)) - \mathrm{E}[f(W)]\}^2] \to 0$ due to $n^{2H-2} \to 0$. I.e., $\frac{1}{n}\sum_{j=1}^n f(U_H(j) - U_H(j-1)) \overset{\mathrm{in}\mathcal{L}^2}{\to} \mathrm{E}[f(W)]$.

Secondly, we apply previous result for (4.11), in fact,

$$\sum_{j=1}^{n} |U_{H}(\frac{j}{n}) - U_{H}(\frac{j-1}{n})|^{\beta}$$

$$= \frac{1}{n^{\beta H}} \sum_{j}^{n} |U_{H}(j) - U_{H}(j-1)|^{\beta}$$

$$= \frac{1}{n^{\beta H-1}} \frac{1}{n} \sum_{j}^{n} |U_{H}(j) - U_{H}(j-1)|^{\beta}$$

$$\stackrel{\text{in L}^{2}}{\longrightarrow} n^{1-\beta H} \mathbf{E}[|W|^{\beta}]$$

Due to $E[|W|^{\beta}] < \infty$, (4.11) holds as well as $n \to \infty$.

THEOREM 4.23. fBm is not a semimartingale for $H \neq \frac{1}{2}$.

Proof. Without loss of generality, we set the time scale T = [0, 1]. Choosing $\beta = 2$ in (4.11), we suppose $U_H(t)$ were a semimartingale.

Case $H < \frac{1}{2}$. Then $\sum_{j=1}^{n} |U_H(\frac{j}{n}) - U_H(\frac{j-1}{n})|^2 \to \infty$ contradicts that semimartingale has finite quadratic variation.

Case $H > \frac{1}{2}$. $\sum_{j=1}^{n} |U_H(\frac{j}{n}) - U_H(\frac{j-1}{n})|^2 \to 0$. On the one hand, according to Doob-Meyer decomposition, $U_H(t) = M(t) + A(t)$, where M(t) is a local martingale and A(t) is local finite variation process. Consider now,

$$\sum_{j=1}^{n} |A(\frac{j}{n}) - A(\frac{j-1}{n})|^{2}$$

$$\leq \sum_{j=1}^{n} |A(\frac{j}{n}) - A(\frac{j-1}{n})| \cdot \sup_{j} |A(\frac{j}{n}) - A(\frac{j-1}{n})|$$

$$\xrightarrow{n\uparrow\infty} \quad 0$$

Hence A(t) has quadratic variation zero and we have $0 = [U_H, U_H] = [M, M]$, where $[\cdot, \cdot]$ is denoted for quadratic variation. Consequently, M(t) is zero process due to Cauchy Schwarz inequality. In other words, $U_H(t) = A(t)$ which has finite variation. On the other Hand, choosing $1 < \gamma < \frac{1}{H}$, then $\sum_{j=1}^{n} |U_H(\frac{j}{n}) - U_H(\frac{j-1}{n})|^{\gamma} \to \infty$. Precisily,

$$\infty \leftarrow \sum_{j=1}^{n} |U_{H}(\frac{j}{n}) - U_{H}(\frac{j-1}{n})|^{\gamma}
\leq \underbrace{\sup_{1 \leq j \leq n} |U_{H}(\frac{j}{n}) - U_{H}(\frac{j-1}{n})|^{\gamma-1}}_{(\gamma-1) \xrightarrow{\text{H\"older}} 0} \cdot \sum_{j=1}^{n} |U_{H}(\frac{j}{n}) - U_{H}(\frac{j-1}{n})|,$$

4.4 fBm is not Semimartingale for $H \neq \frac{1}{2}$

this leads to $\sum_{j=1}^{n} |U_H(\frac{j}{n}) - U_H(\frac{j-1}{n})| \to \infty$, which contradicts mentioned above that $U_H(t)$ has finite variation.

Given all that, fBm is not a semimartingale for $H \neq \frac{1}{2}$.

5 Fractional Ornstein-Uhlenbeck Process

In this section we turn our attention on the fractional Ornstein-Uhlenbeck process (for short, we denote it by fOU).

5.1 Fractional Ornstein-Uhlenbeck Process

Consider the following stochastic dynamics

$$dX_t = -aX_t dt + \gamma dU_H(t), \tag{5.1}$$

where $(X_t)_{t\geq 0}$ is a stochastic process, $a, \gamma \in \mathbb{R}_+$ and $(U_H(t))_{t\geq 0}$ fBm with Hurst exponent H. In fact, given an initial condition $X_0(\omega) = b(\omega)$, then in the theory of SDE, (5.1) is understood as

$$X_t(\omega) = b(\omega) - a \int_0^t X_u(\omega) du + \gamma U_H(t)(\omega)$$
 (5.2)

for t > 0.

DEFINITION 5.1. The *fractional Ornstein-Uhlenbeck* (fOU) process is defined as the solution of (5.1).

If $H = \frac{1}{2}$, fOU is said to be *Ornstein-Uhlenbeck process*. Cheridito et al.[4] shows following Lemmas with Hölder continuity of $U_H(t)$.

LEMMA 5.2. Let $U_H(t)$ be a fBm, $a \in \mathbb{R}_+$ and $s, d \in \mathbb{R}$ such that $d \leq s$. Then there exists a Riemann-Stieljes integral such that

$$\int_{d}^{s} e^{au} dU_{H}(u) = e^{as} U_{H}(s) - e^{ad} U_{H}(d) - a \int_{d}^{s} U_{H}(u) e^{au} du.$$
 (5.3)

Proof. See. [4], p.11, Proposition A.1.

Lemma 5.3. Let $H \in (0, \frac{1}{2}) \cup (\frac{1}{2}, 1), a > 0$ and $-\infty \le m < n \le j < k < \infty$. Then

$$E\left[\int_{m}^{n} e^{au} dU_{H}(u) \int_{j}^{k} e^{as} dU_{H}(s)\right]$$

$$= H(2H - 1) \int_{m}^{n} e^{au} \left(\int_{j}^{k} e^{as} (s - u)^{2H - 2} ds\right) du.$$

Proof. See. [4], p.5, Lemma 2.1.

THEOREM 5.4. $\hat{X}_t^{b,H} := e^{-at} \left(b + \gamma \int_0^t e^{au} dU_H(u) \right)$ is the solution that solves (5.2) for $u \ge 0$.

Proof. Cf. [4], p.11, Proposition A.1. We define

$$Y(t) := \int_0^t X_u \, du,$$

for $t \ge 0$. Rewirte (5.2) with Y_t and Y(0) = 0,

$$Y'(t) = b - aY(t) + \gamma U_H(t)$$

And the solution of that linear differential equation with Y(0) = 0 is

$$Y(t) = e^{-at} \int_0^t e^{au} (b + \gamma U_H(u)) du,$$

in terms of definition above, using (5.3)

$$X(t) = Y'(t)$$

$$= -ae^{-at} \int_{0}^{t} e^{au}(b + \gamma U_{H}(u)) du + e^{-at}e^{at}(b + \gamma U_{H}(t))$$

$$= -ae^{-at} \int_{0}^{t} e^{au}(b + \gamma U_{H}(u)) du + b + \gamma U_{H}(t)$$

$$= e^{-at} \left(-a \int_{0}^{t} e^{au} \gamma U_{H}(u) du + e^{at} \gamma U_{H}(t) - a \int_{0}^{t} e^{au} du \cdot b \right) + b$$

$$= e^{-at} \left(\gamma \int_{0}^{t} e^{au} dU_{H}(u) - e^{au} |_{u=0}^{u=t} \cdot b \right) + b$$

$$= e^{-at} \left(\gamma \int_{0}^{t} e^{au} dU_{H}(u) + b \right)$$

In order to have a stationary solution, we assume that the initial value is centered Gaussian that $\hat{X}_{H,t} := \hat{X}_t^{\gamma \int_{-\infty}^0 e^{au} \, dU_H(u), H} := e^{-at} \left(\gamma \int_{-\infty}^t e^{au} \, dU_H(u) \right).$

THEOREM 5.5. $(\hat{X}_{H,t})_{t\geq 0}$ is centered Gaussian and stationary.

Proof. For the sake of simplicity, we let $\hat{X}_{H,t} = \int_{-\infty}^{t} e^{au} dU_H(u)$. Fix $\epsilon^j > 0, H \in (0,1)$,

then there exists $\{u_0^j < u_1^j < \dots < u_{k_j}^j \le t_j\}$ such that

$$|\hat{X}_{H,t_{j}} - \sum_{l=0}^{k_{j}-1} e^{au_{l}^{j}} (U_{H}(u_{l+1}^{j}) - U_{H}(u_{l}^{j}))|$$

$$= |\int_{-\infty}^{t_{j}} e^{au} dU_{H}(u) - \sum_{l=0}^{k_{j}-1} e^{au_{l}^{j}} (U_{H}(u_{l+1}^{j}) - U_{H}(u_{l}^{j}))|$$

$$< \epsilon^{j}.$$

for $0 \le j \le d$.

On the one hand, we calculate the characteristic function of $(\hat{X}_{H,t_1},\ldots,\hat{X}_{H,t_d})$ approximately with respect to $\epsilon^1,\ldots,\epsilon^d$.

$$\sum_{j=1}^{d} \xi_j \hat{X}_{H,t_j} \approx \sum_{j=1}^{d} \xi_j \left(\sum_{l=0}^{k_j - 1} e^{au_l^j} (U_H(u_{l+1}^j) - U_H(u_l^j)) \right)$$

Notice that since $(U_H(t))$ is centered Gaussian process, any infinite linear combination of its instances is centered Gaussian again. In other words,

$$\sum_{j=1}^{d} \xi_j \left(\sum_{l=0}^{k_j-1} e^{au_l^j} (U_H(u_{l+1}^j) - U_H(u_l^j)) \right)$$

is centered Gaussian. Passing $\epsilon^1, \ldots, \epsilon^d$ to zero, it implies immediately $(\hat{X}_{t_1}, \ldots, \hat{X}_{t_d})$ is centered Gaussian process due to Corolary 2.15.

On the other hand,

$$\mathbb{E}[\exp\{i\sum_{j=1}^{d}\xi_{j}\hat{X}_{H,t_{j}}\}] \approx \mathbb{E}[\exp\{i\sum_{j=1}^{d}\xi_{j}\left(\sum_{l=0}^{k_{j}-1}e^{au_{l}^{j}}(U_{H}(u_{l+1}^{j})-U_{H}(u_{l}^{j}))\right)\}]$$

Since $\{U_H(t)\}\$ has stationary increments, then

$$E[\exp\{i\sum_{j=1}^{d} \xi_{j} \left(\sum_{l=0}^{k_{j}-1} e^{au_{l}^{j}} (U_{H}(u_{l+1}^{j}) - U_{H}(u_{l}^{j}))\right)\}]$$

$$= E[\exp\{i\sum_{j=1}^{d} \xi_{j} \left(\sum_{l=0}^{k_{j}-1} e^{au_{l}^{j}} (U_{H}(u_{l+1}^{j} + \tau) - U_{H}(u_{l}^{j} + \tau))\right)\}]$$

which converges if ϵ 's tend to zero and must be equal $\mathbb{E}[\exp\{i\sum_{j=1}^d \xi_j \hat{X}_{H,t_j+\tau}\}]$, i.e., $(\hat{X}_{H,t})$ is stationary process.

THEOREM 5.6. Let H be that $H \in (0, \frac{1}{2}) \cup (\frac{1}{2}, 1)$. Then

$$\varsigma_{\hat{X}_H}(\tau) = \frac{1}{2} \gamma^2 \sum_{k=1}^N a^{-2k} \left(\prod_{j=0}^{2k-1} (2H - j) \right) \tau^{2(H-k)} + \mathcal{O}(\tau^{2H-2N-2})$$

for $N \in \mathbb{N}, \tau \in \mathbb{R}, \gamma, a \in \mathbb{R}_+$ as in (5.2).

Proof. See [4], p.7, Theorem 2.3.

COROLLARY 5.7. $(\hat{X}_{H,t})_{t\geq 0}$ has long memory for $H\in (\frac{1}{2},1)$.

Proof. Consider the autocovariance of $(\hat{X}_{H,t})$, with a given function $c(a,\gamma,H):=\frac{1}{2}\gamma^2a^{-2k}\left(\prod_{j=0}^{2k-1}(2H-j)\right)$,

$$\varsigma_{\hat{X}_H}(\tau) = \sum_{k=1}^{N} c(a, \gamma, H) \tau^{2H-2k} + \mathcal{O}(\tau^{2H-2N-2})$$

is obviously tending to zero as τ goes to infinity. And in order to check convergence of $\sum_{\tau=1}^{\infty} \varsigma_{\hat{X}_H}(\tau)$, we only need to check the term of k=1 in the summation, because in the case k>1, $\tau^{2H-2k}<\tau^{2H-2}$ for $\tau\in\mathbb{N}$. Note that, $0< c(a,\gamma,H)<\infty$ when $H\in(\frac{1}{2},1)$. We deal with it in the same way as in Theorem 4.18 (use limit comparison test). That means, if 2H<1, namely $H<\frac{1}{2},\sum_{\tau=1}^{\infty}\varsigma_{X_H}(\tau)$ converges. For $H\in(\frac{1}{2},1)$, it diverges. Thus $(\varsigma_{\hat{X}_H}(t))_{t\geq 0}$ has long memory for $H\in(\frac{1}{2},1)$.

6 Applications in Financial Mathematics

6.1 Fractional Black-Scholes Model

In this subsection, we'll introduce fBm to the Black-Scholes model. To be specific, our finance market is modeled with two stochastic processes, i.e. a process of a riskless asset $(A_t)_t$ and a process of price of a stock $(S_t)_t$. The stock is assumed that it pays no dividends. Setting initial conditions $A_0 = 1, S_0 = 1$, we give our fractional Black-Scholes model as follows

$$A_t = \exp(rt)$$

$$S_t = \exp(rt + \mu(t) + \sigma U_H(t)), t \in [0, T], \tag{6.1}$$

where $r \in \mathbb{R}$, $\sigma \in \mathbb{R}_+$, $\sup_{t \in [0,T]} \mu(t) < \infty$. Through out this section we denote by $(\mathcal{F}_t^X)_t$ the filtration of a stochatic process $(X_t)_t$ and give definitions

DEFINITION 6.1. A \mathbb{R}^2 -valued stochastic process $(\xi_t^0, \xi_t^1)_{t \in [0,T]}$ is said to be a *strategy* to (6.1), if $\xi_t^0 \in \mathcal{F}_j^A$ and $\xi_t^1 \in \mathcal{F}_j^S$, for $0 \le j \le t$.

DEFINITION 6.2. A stochastic process $(V_t)_{t \in [0,T]}$ is said to be *value process* with respect to strategy (ξ_t^0, ξ_t^1) , if

$$V_t = \xi_t^0 A_t + \xi_t^1 S_t$$

for $t \in [0, T]$.

DEFINITION 6.3. A stochastic process $(\tilde{V}_t)_{t\in[0,T]}$ is said to be discounted value process of a value process $(V_t)_{t\in[0,T]}$ with respect to (ξ_t^0,ξ_t^1) , if

$$\tilde{V}_t = \frac{V_t}{A_t}$$

for $t \in [0, T]$.

Obviously, $\tilde{V}_t = \xi^0 + \xi^1 \tilde{S}_t$ with $\tilde{S}_t = \exp(\mu(t) + \sigma U_H(t))$.

DEFINITION 6.4. A strategy $(\xi_t^0, \xi_t^1)_{t \in T}$ is said to be *self-financing*, if

$$V_T = V_0 + \sum_{j=1}^m \xi_{s_j}^0 (A_{s_j} - A_{s_{j-1}}) + \xi_{s_j}^1 (S_{s_j} - S_{s_{j-1}}).$$
(6.2)

for $0 = s_0 \le s_1 \le \dots \le s_m = T$.

DEFINITION 6.5. A self-financing strategy $(\xi_t^0, \xi_t^1)_{t \in T}$ is said to be have *arbitrage*, if its discounted process satisfies following conditions

- (i) $\mathcal{P}[\tilde{V}_T \tilde{V}_0] = 1$.
- (ii) $\mathcal{P}[\tilde{V}_T > 0] > 0$.

Cheridito shows that in reality, if there exist a minimal amount of time between two successive transactions, the market (6.1) is arbitagefree. In following we complete the proof in [3] of the claim.

LEMMA 6.6. Let $(X_t)_{t\geq 0}$ be a stochastic process continuous in t. If (X_t) is a modification of the process

$$\left(\int_0^t (t-u)^{H-\frac{1}{2}} dB_u\right)_{t>0}$$

for $(B_t)_{t\geq 0}$ a Brownian motion and $H\in (0,\frac{1}{2})\cup (\frac{1}{2},1)$, then

$$\mathcal{P}[\sup_{t \in [a,b]} X_t \le -c] > 0$$

for $c \ge 0$ and $0 < a \le b$.

Proof. See [3], p.15, Lemma 4.2.

THEOREM 6.7. Let $(S_t)_{t\in[0,T]}$ be a stochastic process such that

$$\tilde{S}_t = \exp\left(\mu(t) + \sigma U_H(t)\right),\tag{6.3}$$

where μ, σ are as in (6.1), $U_H(t)$ is a fBm. If there exist

$$\xi_t^1 = f_0 \mathbb{1}_{\{0\}}(t) + \sum_{k=1}^{n-1} f_k \mathbb{1}_{(\tau_k, \tau_{k+1}]}(t)$$

where $t \in [0, T]$, f_k is family of $\mathcal{F}_k^{U_H}$ -measurable function for $k \in \{1, \dots, n-1\}$. $0 = \tau_1 < \dots < \tau_n = T$ are stopping times with respect to $\mathcal{F}_{\tau_k}^{U_H}$ respectively, with $\tau_{k+1} - \tau_k \ge m$ for some m > 0. If there exists a $k \in \{0, \dots, n-1\}$ such that $\mathcal{P}[f_k \ne 0] > 0$, then

$$\mathcal{P}[(\xi^1 \cdot \tilde{S})_T < 0] > 0,$$

where $(\xi^1 \cdot \tilde{S})_T := \sum_{k=1}^n \xi_{\tau_k}^1 (\tilde{S}_k - \tilde{S}_{k-1}).$

Proof. Cf.[3], p.18, Theorem 4.3. For sake of simplicity, we let $\tilde{S}_t = \exp(U_H(t))$ and $f_0 = 0$. Note that, ξ^1 is predictable. Assume $\mathcal{P}[(\xi^1 \cdot \tilde{S})_T < 0] = 0$, then there exisis

$$l = \min\{j : \mathcal{P}[f_j \neq 0] > 0, \quad \mathcal{P}[\sum_{k=1}^{j} f_k(e^{U_H(\tau_{k+1})} - e^{U_H(\tau_k)}) \ge 0] = 1\}$$

Then either

$$\sum_{k=1}^{j-1} f_k(e^{U_H(\tau_{k+1})} - e^{U_H(\tau_k)}) = 0$$

a.s., or

$$P\left[\sum_{k=1}^{j-1} f_k(e^{U_H(\tau_{k+1})} - e^{U_H(\tau_k)}) < 0\right] < 0.$$

This leads to

$$\mathcal{P}\left[\left(\sum_{k=1}^{l-1} f_k(e^{U_H(\tau_{k+1})} - e^{U_H(\tau_k)}\right) \le 0\right] = 1$$

Ignoring constant term, we define

$$U_H(t)(\omega) = \int_{\mathbb{R}} \mathbb{1}_{\{t \ge u\}} \cdot (t - u)^{H - \frac{1}{2}} - \mathbb{1}_{\{u \le 0\}} (-u)^{H - \frac{1}{2}} d\omega(u).$$

where $\omega(u) := B_u(\omega)$ for all $\omega \in \Omega^w$. We give the filtration $(\mathcal{F}_t^{\Omega^w})$ denoting by

$$\mathcal{F}_t^{\Omega^w} := \sigma(\{\{w \in \Omega^w : \omega(u) \in \mathbb{R}\} : -\infty < u \le t, t \in \mathbb{R}\}).$$

Then τ_k is also stopping time of $\mathcal{F}_t^{\Omega^w}$ due to

$$\mathcal{F}_t^{U_H} \subset \mathcal{F}_t^{\Omega^w}, t \in \mathbb{R}$$

For $\omega \in \Omega^w$, we split it at the time point $\tau_l(\omega)$ into two parts as follows

$$\psi_{\omega}(u) := \omega(u) \mathbb{1}_{(-\infty, \tau_{l}(\omega)]}(u), u \in \mathbb{R}$$

$$\phi_{\omega}(u) := \omega(\tau_{l}(\omega) + u) - \omega(\tau_{l}(\omega)), u \ge 0.$$

Corresponding to each part, we define

$$\Omega^1 := \{ \psi_\omega \in \mathcal{C}(\mathbb{R}) : \omega \in \Omega^w \}$$
$$\Omega^2 := \{ \phi_\omega \in \mathcal{C}([0, \infty)) : \omega \in \Omega^w \}$$

And for the smallst σ -algebra of all subsets, respectively, of Ω^1, Ω^2 are denoted by $\mathscr{B}^1, \mathscr{B}^2$. Notice that

$$\{ \tau_l \le t \} \cap \{ \psi_\omega \in \Omega^w \} = \{ \{ \omega \in \Omega^w : \omega(u) \in \mathbb{R} \} : -\infty < u \le t \}$$
$$\in \mathcal{F}_t^{\Omega^w}.$$

therefore is ψ_{ω} a $\mathcal{F}_{\tau_l}^{\Omega^w}$ - measurable mapping. Moreover, since the strong Markovian property of Brownian motion, ϕ_{ω} is independent of $\mathcal{F}_{\tau_l}^{\Omega^w}$ and it must be a Brownian motion.

Plugging $\psi_{\omega}, \phi_{\omega}$ into Ω , we calculate the value process

$$\left(\sum_{k=1}^{l-1} f_k(e^{U_H(\tau_{k+1})} - e^{U_H(\tau_k)}) + f_l(e^{U_H(\tau_l + m)} - e^{U_H(\tau_l)})\right)(\omega)$$

$$= \sum_{k=1}^{l-1} f_k(e^{U_H(\tau_{k+1})} - e^{U_H(\tau_k)})(\omega) + f_l(e^{U_H(\tau_l + m)} - e^{U_H(\tau_l)})(\omega)$$

$$= J^1(\omega) + f_l(\exp\{\int_{\mathbb{R}} \mathbb{1}_{\{\tau_l \ge u\}}(\tau_l(\omega) + m - u)^{H - \frac{1}{2}} - \mathbb{1}_{\{\tau_l \ge u\}}(\tau_l(\omega) - u)^{H - \frac{1}{2}} d\omega(u)$$

$$+ \int_{\mathbb{R}} \mathbb{1}_{(\tau_l, \tau_l + m)}(\tau_l + m - u)^{H - \frac{1}{2}} d\omega(u)\})$$

$$= J^1(\psi_\omega) + f_l\left(\exp\{\int_{-\infty}^{\tau_l(\psi_\omega)} (\tau_l(\psi_\omega) + m + u)^{H - \frac{1}{2}} - (\tau_l(\psi_\omega) - u)^{H - \frac{1}{2}} d\psi_\omega(u)\}\right)_{=:J^2(\psi_\omega, m)}$$

$$\cdot \underbrace{\left(\exp\{\int_{0}^{m} (m - u)^{H - \frac{1}{2}} d\phi_\omega(u)\} - 1\right)}_{=:J^3(\phi_\omega, m)}$$

$$= J(\psi_\omega, \phi_\omega, m)$$

where J is defined as

$$J(\psi, \phi, t) := J^{1}(\psi) + J^{2}(\psi, t) \cdot J^{3}(\phi, t)$$

for $\psi \in \Omega^1, \phi \in \Omega^2$.

Indeed, for $\psi \in \Omega^1$, $\phi \in \Omega^2$, $J(\psi, \phi, \cdot)$ has continuous path on $(\Omega^1 \times \Omega^2, \mathcal{B}^1 \otimes \mathcal{B}^2)$, then we can define a $\mathcal{B}^1 \otimes \mathcal{B}^2$ - measurable set

$$E:=\{(\psi,\phi): \sup_{m\leq t\leq T}J(\psi,\phi,t)<0\}$$

Note that

$$E[\mathbb{1}_{E}(\psi_{\omega}, \phi_{\omega}) | \psi_{\omega} = \omega_{1}] = \mathcal{P}[\sup_{m \leq t \leq T} J(\omega_{1}, \phi_{\omega}, t) < 0]$$

$$\geq \mathcal{P}[J^{1}(\omega_{1}) + \sup_{m \leq t \leq T} J^{2}(\omega_{1}, t) \cdot \sup_{m \leq t \leq T} J^{3}(\phi_{\omega}, t) < 0]$$

It is clear $J^2(\omega_1)$ is bounded for a fixed ω_1 on the compact set and $J^1(\omega_1)$ also. We set $J^1(\omega_1) = c_1$, $\sup_{m \le t \le T} J^2(\omega_1, t) = c_2$. Thanks to Lemma 6.6, there exist $D \subset \Omega^w$ such that $\mathcal{P}^w[D] > 0$ and

$$\int_{0}^{m} (m-u)^{H-\frac{1}{2}} d\phi_{\omega}(u) < c_{3}$$

for $\omega \in D$. where c_3 could be small enough that $c_2 \cdot e^{c_3 - 1} < 0$. Under assumption at beginning, $c_1 \leq 0$ almost surely. All of this leads to

$$E[\mathbb{1}_{E}(\omega_{1}, \phi_{\omega}) | \psi_{\omega} = \omega_{1}]$$

$$\geq \mathcal{P}[J^{1}(\omega_{1}) + \sup_{m \leq t \leq T} J^{2}(\omega_{1}, t) \cdot \sup_{m \leq t \leq T} J^{3}(\phi_{\omega}, t) < 0]$$

$$> 0$$

for $\omega \in \Omega^w$. Then

$$\mathcal{P}\left[\sum_{k=1}^{l} f_{k}(e^{U_{H}(\tau_{k+1})} - e^{U_{H}(\tau_{k})}) < 0\right]$$

$$\geq \mathcal{P}\left[\sum_{k=1}^{l-1} f_{k}(e^{U_{H}(\tau_{k+1})} - e^{U_{H}(\tau_{k})}) + \sup_{m \leq t \leq T} f_{l}(e^{U_{H}(\tau_{l}+t)} - e^{U_{H}(\tau_{l})}) < 0\right]$$

$$= \mathbb{E}\left[\mathbb{E}\left[\mathbb{1}_{E}(\omega_{1}, \phi_{\omega}) | \psi_{\omega} = \omega_{1}\right]\right] > 0$$

this contradicts our assumption. It must be that $\mathcal{P}[(\xi^1 \cdot S)_T < 0] > 0$.

COROLLARY 6.8. Let strategy (ξ^0, ξ^1) be such that, ξ^1 is given as in Theorem 6.7. Then the strategy has no arbitrage in our finance market (6.1).

Proof. Assume (ξ^0, ξ^1) is a self-financing strategy. In terms of Definition 6.4,

$$\tilde{V}_T - \tilde{V}_0 = \sum_{k=1}^n \frac{(\xi_k^0 A_k + \xi_k^1 S_k)}{A_k} - \frac{(\xi_k^0 A_{k-1} + \xi_k^1 S_{k-1})}{A_{k-1}}$$

$$= \sum_{k=1}^n \xi_k^1 (\tilde{S}_k - \tilde{S}_{k-1})$$

It follows then from Theorem 6.7, $\mathcal{P}[(\tilde{V}_T - \tilde{V}_0) < 0] > 0$, since and therefore the strategy has no arbitrage in (6.1).

6.2 Fractional Calculus and Discretization of Fractionally Integrated Process

The fractional integral could be derived from the repeated integral which is approached by Riemann-Liouville integral.

$$\int_0^s \int_0^{s_1} \int_0^{s_2} \cdots \int_0^{s_{n-1}} f(s_n) ds_n \cdots ds_2 ds_1$$

$$= \frac{1}{(n-1)!} \int_0^s (s-u)^{n-1} f(u) du$$

for $n \in \mathbb{N}$. Where n is said to be the order of the fractional integral. We extend it with the order $\alpha \in \mathbb{R}_+$.

DEFINITION 6.9. Let f be a locally integrable function. The Riemann-Liouville fractional integral of order α is defined as follows

$$I^{\alpha}f(s) := \frac{1}{\Gamma(\alpha)} \int_{0}^{s} (s-u)^{\alpha-1} f(u) \, du \tag{6.4}$$

for $s, \alpha \in \mathbb{R}_+$.

Remark, $I^{\alpha}I^{\beta}=I^{\beta}I^{\alpha}=I^{\alpha+\beta}$. Let $\Phi_{\alpha}(s):=\frac{s^{\alpha-1}}{\Gamma(\alpha)}$, then

$$I^{\alpha}f(s) = \Phi_{\alpha}(s) * f(s) \tag{6.5}$$

We give the definition of fractional derivative.

DEFINITION 6.10. Let f be a locally integrable function. The Riemann-Liouville fractional derivative of order α is defined as follows

$$D^{\alpha}f(s) := \begin{cases} \frac{d^n}{ds^n} \left[\frac{1}{\Gamma(n-\alpha)} \int_0^s \frac{f(u)}{(s-u)^{\alpha+1-n}} du \right] & \text{if } n-1 < \alpha < n, \\ \frac{d^n}{ds^n} f(s) & \text{if } \alpha = n \end{cases}$$
(6.6)

for $n \in \mathbb{Z}$ and $\alpha \in \mathbb{R}, s \in \mathbb{R}_+$. In [5], Comte gave a truncated representation fBm $(U_{\alpha,t})$.

$$U_{\alpha,t} = \int_0^t \frac{(t-u)^{\alpha}}{\Gamma(1+\alpha)} dB_u, \tag{6.7}$$

Where $|\alpha| < \frac{1}{2}, (B_u)$ Brownian motion. Notice that α is given by $H - \frac{1}{2}$ as in previous representation and the long memory property remains true.

DEFINITION 6.11. The fractionally integrated process of order α , $|\alpha| < \frac{1}{2}$ is defined as

$$X(t) = \int_0^t g(t-u) \, dB_u \tag{6.8}$$

where $(B_t)_t$ is Brownian motion and

$$g(s) = \Phi_{\alpha+1}(s)h(s) \tag{6.9}$$

$$g(s) = \Phi_{\alpha+1}(s)h(s)$$

$$= \frac{s^{\alpha}h(s)}{\Gamma(1+\alpha)}$$
(6.9)

with $h \in C^1([0, T])$.

PROPOSITION 6.12. Let X(t) be a fractionally integrated process of order α , $|\alpha| < \frac{1}{2}$, then

$$X(t) = \int_{0}^{t} c(t-u) \, dU_{\alpha,u} \tag{6.11}$$

$$:= \frac{d}{dt} \left(\int_0^t c(t-u) U_{\alpha,u} \, du \right) \tag{6.12}$$

where $c \in C([0,T])$, c and g are functions related by

$$c(s) = \frac{d}{ds} \left(\int_0^s \frac{(s-u)^{-\alpha} u^{\alpha} h(u) du}{\Gamma(1-\alpha)\Gamma(1+\alpha)} \right). \tag{6.13}$$

$$u^{\alpha}h(u) = \frac{d}{du}\left(\int_0^u c(s)(u-s)^{\alpha} ds\right). \tag{6.14}$$

Proof. See [6], pp.106-108, Lemma 1, Proposition 1.

DEFINITION 6.13. Let X_t be a fractionally integrated process of order α on [0,T] and $|\alpha| < \frac{1}{2}$. The fractional derivation of order α is defined as

$$X^{(\alpha)}(t)$$

$$:= \int_0^t \frac{(t-u)^{-\alpha}}{\Gamma(1-\alpha)} dX_t$$

$$:= \frac{d}{dt} \int_0^t \frac{(t-u)^{-\alpha}}{\Gamma(1-\alpha)} X_t du$$

PROPOSITION 6.14. $X^{(\alpha)}(t)$ is well-defined and mean square continuous. If h(0) is invertible and $h \in C^2([0,T])$, then $X^{(\alpha)}(t)$ has the $MA(\infty)$ representation

$$X^{(\alpha)}(t) = \int_0^t c(t-s) dB_s.$$

where c and h are one-to-one related by (6.13) and (6.14).

Proof. See [6], p.111, Proposition 4.

THEOREM 6.15. Let X_t be a locally integrable function on [0,T] and $|\alpha|<\frac{1}{2}$. Then

$$X(t) = \int_0^t \frac{(t-u)^{\alpha}}{\Gamma(1+\alpha)} dX_u^{(\alpha)}$$
(6.15)

Proof. Since all the integrands is nonnegative, we could apply Fubini theorem.

$$\int_0^t \frac{(t-u)^{\alpha}}{\Gamma(1+\alpha)} dX_u^{(\alpha)}$$

$$= \frac{d}{dt} \int_0^t \frac{(t-u)^{\alpha}}{\Gamma(1+\alpha)} X_u^{(\alpha)} du$$

$$= \frac{d}{dt} \int_0^t \frac{(t-u)^{\alpha}}{\Gamma(1+\alpha)} \left(\int_0^u \frac{(u-s)^{-\alpha}}{\Gamma(1-\alpha)} dX_s \right) du$$

$$= \frac{d}{dt} \int_0^t \left(\int_s^t \frac{(t-u)^{\alpha}(u-s)^{-\alpha}}{\Gamma(1+\alpha)\Gamma(1-\alpha)} du \right) dX_s$$

Changing variable with $v := \frac{u-s}{t-s}$,

$$\int_{s}^{t} \frac{(t-u)^{\alpha}(u-s)^{-\alpha}}{\Gamma(1+\alpha)\Gamma(1-\alpha)} du$$

$$= \int_{0}^{1} \frac{(t-s-(t-s)v)^{\alpha}(t-s)^{-\alpha}}{\Gamma(1+\alpha)\Gamma(1-\alpha)} (t-s) dv$$

$$= \frac{(t-s)}{\Gamma(1+\alpha)\Gamma(1-\alpha)} \int_{0}^{1} (1-v)^{\alpha} v^{-\alpha} dv$$

Note that the beta function

$$B(x,y) = \int_0^1 t^{x-1} (1-t)^{y-1} dt$$
$$= \frac{\Gamma(x)\Gamma(y)}{\Gamma(x+y)}$$

for $x, y \in \mathbb{R}_+$. Then,

$$\int_{s}^{t} \frac{(t-u)^{\alpha}(u-s)^{-\alpha}}{\Gamma(1+\alpha)\Gamma(1-\alpha)} du$$

$$= \frac{(t-s)}{\Gamma(1+\alpha)\Gamma(1-\alpha)} B(1+\alpha, 1-\alpha)$$

$$= \frac{(t-s)}{\Gamma(1+\alpha)\Gamma(1-\alpha)} \frac{\Gamma(1+\alpha)\Gamma(1-\alpha)}{\Gamma(2)}$$

$$= t-s$$

Plugging it back,

$$\int_0^t \frac{(t-u)^\alpha}{\Gamma(1+\alpha)} dX_u^{(\alpha)}$$

$$= \frac{d}{dt} \int_0^t t - s dX_s$$

$$= \frac{d^2}{dt^2} \int_0^t (t-s) X_s ds$$

$$= \frac{d^2}{dt^2} \int_0^t X_s \int_s^t du ds$$

$$= \frac{d^2}{dt^2} \int_0^t \int_s^t X_s du ds$$

$$= \frac{d^2}{dt^2} \int_0^t \int_0^u X_s ds du$$

$$= X_t$$

EXAMPLE 6.16. The truncated fOU process driven by $U_{\alpha,t}$ is

$$F_{\alpha,t} = \gamma \int_0^t e^{-a(t-u)} dU_{\alpha,u} \tag{6.16}$$

In terms of (6.8)

$$c(u) = \gamma e^{-au}$$

And according to (6.14),

$$= \frac{g(u)}{\frac{d}{du} \int_0^u c(s)(u-s)^\alpha ds}{\Gamma(1+\alpha)}$$
$$= \frac{\gamma}{\Gamma(1+\alpha)} \frac{d}{du} \int_0^u e^{-as} (u-s)^\alpha ds$$

Using partial integration,

$$\begin{split} &g(u)\\ &= \frac{\gamma}{\Gamma(1+\alpha)} \left(\left(\frac{d}{du} (-\frac{1}{1+\alpha} e^{-as} (u-s)^{1+\alpha}|_{s=0}^{u}) \right) - \left(\frac{d}{du} (\int_{0}^{u} -ae^{-as} (u-s)^{\alpha}) \, ds \right) \right) \\ &= \frac{\gamma}{\Gamma(1+\alpha)} \left(\left(\frac{d}{du} (\frac{u^{1+\alpha}}{1+\alpha}) \right) - ae^{-au} \left(\int_{0}^{u} e^{-as} s^{\alpha} \, ds \right) \right) \\ &= \frac{\gamma}{\Gamma(1+\alpha)} (u^{\alpha} - ae^{-au} \int_{0}^{u} e^{-as} s^{\alpha} \, ds) \end{split}$$

According to (6.10), $h(u) = \gamma(1 + \frac{ae^{-au} \int_0^u e^{-as} s^{\alpha} ds}{u^{\alpha}})$. It is clear $h'(0) \neq 0$ and $h \in C^2([0,T])$. I.e., its fractional derivative of order α is

$$F_t^{(\alpha)} = \int_0^t c(t-u) dB_u$$
$$= \gamma \int_0^t e^{-a(t-u)} dB_u$$
(6.17)

and

$$F_{\alpha,t} = \int_0^t \frac{(t-s)^{\alpha}}{\Gamma(1+\alpha)} dF^{(\alpha)}(s). \qquad (6.18)$$

Obviously, $F_t^{(\alpha)}$ is the solution of follows (as by fOU, c.f. Theorem 5.4)

$$dF^{(\alpha)}(t) = -aF^{(\alpha)}(t) dt + \gamma dB_t, \quad F^{(\alpha)}(0) = 0$$

I.e., $(F_t^{(\alpha)})_{t\geq 0}$ is the Ornstein-Uhlenbeck process.

In order to approximate a fOU $F_{\alpha,t}$ on a discrete time scale, according to (6.18), Comte and Renault define an approximation by step functions

$$\tilde{F}_{\alpha,n}(t) = \sum_{k=1}^{\lfloor nt \rfloor} \frac{(t - \frac{k-1}{n})^{\alpha}}{\Gamma(1+\alpha)} \left(F_{\frac{k}{n}}^{(\alpha)} - F_{\frac{k-1}{n}}^{(\alpha)} \right) \tag{6.19}$$

for $n \in \mathbb{N}, t \in \mathbb{R}_+$.

THEOREM 6.17. $\tilde{F}_{\alpha,n}(t) \to F_{\alpha,n}(t)$ in distribution as n tends to infinity.

Proof. See. [7], p.298, Proposition 3.1.

Notice that, $F^{(\alpha)}(t)$ is an AR(1) process due to (6.17), i.e.

$$(1 - c_n L_n) F_{\frac{k}{n}}^{(\alpha)} = \epsilon(\frac{k}{n}) \tag{6.20}$$

where $\epsilon(\cdot) \sim \mathcal{N}(0, \operatorname{Var}[F_{\cdot}^{(\alpha)}]), L_n$ represents the lag operator such that $L_n F_{\frac{k}{n}}^{(\alpha)} = F_{\frac{k-1}{n}}^{(\alpha)}$ and c_n are normalized coefficients corresponding to (6.19). Moreover, on the one hand, as n goes to infinity, we have asymptotically following equation in distribution.

$$F_{\alpha,n}(\frac{k}{n}) = \sum_{l=1}^{k} \frac{(\frac{k}{n} - \frac{l-1}{n})^{\alpha}}{\Gamma(1+\alpha)} \left(F_{\frac{l}{n}}^{(\alpha)} - F_{\frac{l-1}{n}}^{(\alpha)} \right)$$

$$= \sum_{l=1}^{k} \frac{(k-l+1)^{\alpha}}{\Gamma(1+\alpha)n^{\alpha}} \left(F_{\frac{l}{n}}^{(\alpha)} - F_{\frac{l-1}{n}}^{(\alpha)} \right)$$

$$\stackrel{m=k-l}{=} \sum_{m=0}^{k-1} \frac{(m+1)^{\alpha}}{\Gamma(1+\alpha)n^{\alpha}} \left(F_{\frac{k-m}{n}}^{(\alpha)} - F_{\frac{k-m-1}{n}}^{(\alpha)} \right)$$

$$= \left(\sum_{m=0}^{k-1} \frac{(m+1)^{\alpha} - m^{\alpha}}{\Gamma(1+\alpha)n^{\alpha}} \cdot L_{n}^{m} \right) F_{\frac{k}{n}}^{(\alpha)}$$

$$\stackrel{(6.20)}{=} \left(\sum_{m=0}^{k-1} \frac{(m+1)^{\alpha} - m^{\alpha}}{\Gamma(1+\alpha)n^{\alpha}} \cdot L_{n}^{m} \right) (1-c_{n}L_{n})^{-1} \epsilon(\frac{k}{n}). \tag{6.21}$$

We rewrite this then

$$(1 - c_n L_n) \left(\sum_{m=0}^{k-1} \frac{(m+1)^{\alpha} - m^{\alpha}}{\Gamma(1+\alpha)n^{\alpha}} \cdot L_n^m \right)^{-1} F_{\alpha,n}(\frac{k}{n}) = \epsilon(\frac{k}{n}).$$

On the other hand, if we take consideration with the $ARFIMA(1, \alpha, 0)$, replaced the content with $\{\}^{-1}$ by the α -integrated term $(1 - L_n)^{\alpha}$. Then we have

$$(1 - c_n L_n)(1 - L_n)^{\alpha} F_{\alpha,n}(\frac{k}{n}),$$

which is in general not Gaussian. Furthermore, $(1 - L_n)^{\alpha} F_{\alpha,n}$ is in general not a AR(1). In other words, the ARFIMA(1, α , 0) may not fit such a high-frequency data that is comfortable with the fractional stochastic volatility model driven by fBm.

6.3 Fractional Stochastic Volatility Model

In framework of the Black-Scholes model, a risky asset price is modelded as follows

$$dS_t = r_t S_t dt + \sigma_t S_t dB_t, \quad t \in [0, T], \tag{6.22}$$

for $r, \sigma \in \mathbb{R}_+$. In the simplest case, the volatility is assumed as a constant or a deterministic function of time and underlying price of the asset. Such models, however, generate unrealistic volatility dynamics. To solve this problem, in such Hull-White, Heston or SABR models, the volatility also modelded as a stochastic process(e.g. semimartingale).

In this thesis, the log-volatility $\log(\sigma)_{t\geq 0}$ is assumed to obey fractional Ornstein-Uhlenbeck process $(X_t)_{t\geq 0}$. That means

$$\sigma_t = \exp\{X_t\}$$

$$dX_t = -aX_t dt + \gamma dU_H(t)$$
(6.23)

where $a, \gamma \in \mathbb{R}_+$. In the proceeding section, we have a stationary solution

$$\hat{X}_{H,t} = e^{-at} \int_{-\infty}^{t} e^{au} \, dU_H(u) \tag{6.24}$$

for an appropriate initial condition. Note that, recall the (6.16), $\hat{X}_{H,t} - F_{H-\frac{1}{2},t} = e^{-at}(X_0 - F_{H-\frac{1}{2},0}) \to 0$, as $t \to \infty$. We set $\hat{\sigma}_{H,t} := \exp{\{\hat{X}_{H,t}\}}$, which has following property.

PROPOSITION 6.18. Let $\hat{X}_{H,t}$ be such that as in (6.24) and $\hat{\sigma}_{H,t} = \exp{\{\hat{X}_{H,t}\}}$, then $(\hat{\sigma}_{H,t})$ is weak stationary and has long memory for $H \in (\frac{1}{2}, 1)$.

Proof. We start our proof by definition of the autocovariance of $\hat{\sigma}_{H,t}$

$$\begin{split} &\varsigma_{\hat{\sigma}_{H}}(\tau) \\ &= \operatorname{Cov}[\hat{\sigma}_{H,t}, \hat{\sigma}_{H,t+\tau}] \\ &= \operatorname{E}[\hat{\sigma}_{H,t}\hat{\sigma}_{H,t+\tau}] - \operatorname{E}[\hat{\sigma}_{H,t}] \operatorname{E}[\hat{\sigma}_{H,t+\tau}] \\ &= \operatorname{E}[\exp(\hat{X}_{H,t} + \hat{X}_{H,t+\tau})] - \operatorname{E}[\exp(\hat{X}_{H,t})] \operatorname{E}[\exp(\hat{X}_{H,t+\tau})] \end{split}$$

Since $\hat{X}_{H,t}$, $\hat{X}_{H,t+\tau}$ are centred Gaussian, we apply (2.7) for it

$$\varsigma_{\hat{\sigma}_{H}}(\tau) = \exp(\frac{1}{2} \operatorname{Var}[\hat{X}_{H,t} + \hat{X}_{H,t+\tau}]) - \exp(\frac{1}{2} \operatorname{Var}[\hat{X}_{H,t}]) \exp(\frac{1}{2} \operatorname{Var}[\hat{X}_{H,t+\tau}])$$

Since $(\hat{X}_{H,t})_t$ is stationary, we have

$$\varsigma_{\hat{\sigma}_{H}}(\tau)
= \exp(\frac{1}{2} \text{Var}[\hat{X}_{H,t} + \hat{X}_{H,t+\tau}]) - \exp(\text{Var}[\hat{X}_{H,t}])
= \exp(\frac{1}{2} (\text{Var}[\hat{X}_{H,t}] + \text{Var}[\hat{X}_{H,t+\tau}] + 2\text{E}[\hat{X}_{H,t}\hat{X}_{H,t+\tau}])) - \exp(\text{Var}[\hat{X}_{H,t}])
= \exp(\text{Var}[\hat{X}_{H,t}] + \text{Cov}[\hat{X}_{H,t}, \hat{X}_{H,t+\tau}]) - \exp(\text{Var}[\hat{X}_{H,t}])$$

The term $\operatorname{Var}[\hat{X}_{H,t}]$ is independent of τ . We define $\operatorname{Var}[\hat{X}_{H,t}] = C$. Then

$$\varsigma_{\hat{\sigma}_{H}}(\tau)$$

$$= \exp(C) \exp(\varsigma_{\hat{X}_{H}}(\tau)) - \exp(C)$$

$$= \exp(C) (\exp(\varsigma_{\hat{X}_{H}}(\tau)) - 1)$$

$$= \kappa(\exp(\varsigma_{\hat{X}_{H}}(\tau)) - 1)$$

for some κ . Obviously, $\mathbb{E}[\hat{\sigma}_H(t)] = 1$. With $\varsigma_{\hat{\sigma}_H}(0) = \kappa(\exp(\varsigma_{\hat{X}_H}(0)) - 1)$, the first claim is proved.

On the one hand, consider in Theorem 5.6,

$$\varsigma_{\hat{X}_H}(\tau)) = \mu \tau^{2H-2} + \mathcal{O}(\tau^{2H-2N-2})$$

for some μ . $\zeta_{\hat{X}_H}(\tau)$ vanishes, for $H \in (\frac{1}{2}, 1)$, as τ goes to infinity. I.e.,

$$\lim_{\tau \to \infty} \kappa(\exp(\varsigma_{\hat{X}_H}(\tau)) - 1) = 0.$$

On the other hand, And $\zeta_{\hat{X}_H}(\tau)$ is the equivalence infinitesimal of $\exp(\zeta_{\hat{X}_H}(\tau)) - 1$. Hence,

$$\lim_{\tau \to \infty} \left| \frac{\varsigma_{\hat{\sigma}_{H}}(\tau) - \kappa \mu \tau^{2H-2}}{\tau^{2H-2N-2}} \right|$$

$$= \lim_{\tau \to \infty} \left| \frac{\kappa(\exp(\varsigma_{\hat{X}_{H}}(\tau)) - 1) - \kappa \mu \tau^{2H-2}}{\tau^{2H-2N-2}} \right|$$

$$= \lim_{\tau \to \infty} \left| \frac{\kappa(\varsigma_{\hat{X}_{H}}(\tau) - \mu \tau^{2H-2})}{\tau^{2H-2N-2}} \right|$$

$$< \infty$$

This implies $\varsigma_{\hat{\sigma}_H}(\tau) = \kappa \mu \tau^{2H-2} + \mathcal{O}(\tau^{2H-2N-2})$. For the same reason as in Theorem 4.18, the long memory property requires 1 < 2H, namely, $H \in (\frac{1}{2}, 1)$.

The long memory property may explain why in stock market, large upheavals tend to be followed by large upheavals and small upheavals often happen after by small upheavals. In order to model long memory volatility process, Comte and Renault are forced to set $H \in (\frac{1}{2}, 1)$ in (6.24) named fractional volatility stochastic (FSV), see[7].

6.4 Rough Fractional Stochastic Volatility Model

In order to generate a desirable volatility dynamics, Gatheral et al.(2014) take the implied volatility $\sigma^{BS}(m,\tau)$ into account, where m is the log-moneyness and τ is time to expiration date. The implied volatility refer to the value of volatility required in the Black-Scholes

model such that the pricing is coincide with the asset price we observed. Graphing implied Volatility as a function of moneyness and time to expiration seems to be a U-shape which is so-called volatility smile. In particular, the term structure of volatility skew of at-themoney of stylized data

$$\kappa(\tau) = \left| \frac{\partial}{\partial m} \sigma^{BS}(m, \tau) \right|_{m=0},$$

acts as a power law with exponent around $-\frac{1}{2}$, which is explained for example in [10]. On the one hand, in the FSV model with $H \in (\frac{1}{2}, 1)$, the volatility smile is depressed by arising τ , see [8], p.350, Eq. (4.7). On the other hand, in [9], the volatility is driven by FBM with Hurst exponent H in Fukasawa's model whose volatility skew of at-the-money has a form $\kappa(\tau) \sim \tau^{H-\frac{1}{2}}$ as τ goes to zero. This requires that H is near zero to match the power law decay of $\kappa(\tau)$. All of this suggest us, we should apply a stochastic volatility model which is driven by FBM with Hurst exponent $H \in (0, \frac{1}{2})$.

Replaced by $H \in (0, \frac{1}{2})$ in (6.23), a 'rough fractional stochastic volatility model' (RFSV) is

$$dX_t = -aX_t dt + \gamma dU_H(t), \quad t \in [0, T]$$

$$(6.25)$$

for $a, \gamma \in \mathbb{R}_+$ and a stationary solution

$$\hat{X}_{H,t} = e^{-at} \gamma \int_{-\infty}^{t} e^{au} dU_H(u). \qquad (6.26)$$

Consider a quantity defined on [0,T] with mesh τ

$$s(\tau, \sigma) = \frac{1}{N} \sum_{k=1}^{N} |\log(\sigma_{k\tau}) - \log(\sigma_{(k-1)\tau})|^2$$
(6.27)

Where $N = \lfloor T/\tau \rfloor$. This quantity could describe the smoothness of $(\sigma_t)_t$. Due to the volatilities are not observable, we could take spot volatility values to estimate them. In [3], the daily spot variances by daily realized variance estimates $\tilde{\sigma}$ are used. A plotting of $\log(s(\tau,\tilde{\sigma}))$ against $\log(\tau,\tilde{\sigma})$ looks then as straight line, i.e.,

$$s(\tau, \tilde{\sigma}) = k\tau^z. \tag{6.28}$$

RFSV model does match the observed phenomenon. Replaced $\sigma(t)$ by $\hat{X}_{H,t}$

$$s(\tau, \hat{X}_H) = \frac{1}{N} \sum_{k=1}^{N} |\hat{X}_{H,k\tau} - \hat{X}_{H,(k-1)\tau}|^2$$

Since $(\hat{X}_{H,t})$ is stationary, we could apply weak law of large number, as N goes to infinity,

$$s(\tau, \hat{X}_H) \stackrel{\tau\downarrow 0}{\to} \text{E}[|\hat{X}_{H,t+\tau} - \hat{X}_{H,t}|^2]$$

$$= 2\text{Var}[\hat{X}_{H,t}] - 2\text{Cov}[\hat{X}_{H,t}, \hat{X}_{H,t+\tau}]$$

THEOREM 6.19. Let $(\hat{X}_{H,t,a})$ be as in (6.26) driven by a fBm $(U_H(t))$ with $H \in (0, \frac{1}{2})$, then

$$\mathbb{E}[\sup_{t\in[0,T]}|\hat{X}_{H,t+\tau,a}-\gamma U_H(t)|]\to 0$$

as a goes to zero for T > 0.

Proof. Cf. [10], p.15, Proposition
$$3.1$$
.

The theorem shows if a is small enough, $(\hat{X}_{H,t})_t$ behaves essentially as fBm at a compact time scale.

THEOREM 6.20. Let $(\hat{X}_{H,t,a})$ the solution by (6.26) with $H \in (0,\frac{1}{2})$, then

$$\operatorname{Var}[\hat{X}_{H,t,a}] - \operatorname{Cov}[\hat{X}_{H,t,a}, \hat{X}_{H,t+\tau,a}] \to \frac{1}{2} \gamma^2 \tau^{2H}$$
 (6.29)

as a goes to zero, for $t > 0, \tau > 0$.

(6.29) shows, the choise of $H \in (0, \frac{1}{2})$ enable us to model the log-volatility process has a form of (6.28). It may turn out to be RFSV is more reasonable than FSV for this empirical result.

6.5 Weighted Fractional Stochastic Volatility Model

DEFINITION 6.21. A mixed fractional Brownian motion is defined as follows

$$M_{\alpha,\beta,H_1,H_2}(t) = \alpha U_{H_1}(t) + \beta U_{H_2}(t) \tag{6.30}$$

for $t \in \mathbb{R}$, where α, β are real numbers and U_{H_1}, U_{H_2} are two independent fBm's with Hurst exponents $H_1 \in (0, \frac{1}{2}), H_2 \in (\frac{1}{2}, 1)$ respectively.

PROPOSITION 6.22. The mixed fractional Brownian motion $M_{\alpha,\beta,H_1,H_2}(t)_{t\in\mathbb{R}}$ has following properties

- (i) $M_{\alpha,\beta,H_1,H_2}(0) = 0$ and $(M_{\alpha,\beta,H_1,H_2}(t))_t$ is a centered Gaussian process.
- (ii) $\operatorname{Cov}[M_{\alpha,\beta,H_1,H_2}(t), M_{\alpha,\beta,H_1,H_2}(s)] = \alpha^2 \operatorname{Cov}[U_{H_1}(t), U_{H_1}(s)] + \beta^2 \operatorname{Cov}[U_{H_2}(t), U_{H_2}(s)] = \frac{1}{2} \left(\alpha^2 (t^{2H_1} + s^{2H_1} + |t s|^{2H_1}) + \beta^2 (t^{2H_2} + s^{2H_2} + |t s|^{2H_2}) \right).$
- (iii) $M_{\alpha,\beta,H_1,H_2}(qt) \sim M_{\alpha q^{H_1},\beta q^{H_2},H_1,H_2}(t)$, for $q \in \mathbb{R}$.

Proof. (i): $E[M_{\alpha,\beta,H_1,H_2}(t)] = \alpha E[U_{H_1}(t)] + \beta E[U_{H_2}(t)] = \alpha \cdot 0 + \beta \cdot 0 = 0$. Consider,

$$\sum_{k=1}^{d} c_k M_{\alpha,\beta,H_1,H_2}(k)$$

$$= \sum_{k=1}^{d} c_k (\alpha U_{H_1}(k) + \beta U_{H_2}(k))$$

$$= \sum_{k=1}^{d} c_k \alpha U_{H_1}(k) + \sum_{k=1}^{d} c_k \beta U_{H_1}(k)$$

Since $U_{H_1}(t), U_{H_2}(t)$ are independent and $(U_{H_1}(t))_t, (U_{H_2}(t))_t$ are centered Gaussian process, $\sum_{k=1}^d c_k \alpha U_{H_1}(k) + \sum_{k=1}^d c_k \beta U_{H_2}(k)$ is centered Gaussian and therefore $(M_{\alpha,\beta,H_1,H_2}(t))_t$ is centered Gaussian process.

(ii): Using independence of $U_{H_1}(t)$ and $U_{H_2}(t)$, we have

$$Cov[M_{\alpha,\beta,H_{1},H_{2}}(t), M_{\alpha,\beta,H_{1},H_{2}}(s)]$$

$$= Cov[\alpha U_{H_{1}}(t) + \beta U_{H_{2}}(t), \alpha U_{H_{1}}(t) + \beta U_{H_{2}}(s)]$$

$$= E[(\alpha U_{H_{1}}(t) + \beta U_{H_{2}}(t))(\alpha U_{H_{1}}(t) + \beta U_{H_{2}}(s))]$$

$$= E[\alpha^{2}U_{H_{1}}(t)U_{H_{1}}(s)] + \underbrace{E[\alpha\beta U_{H_{2}}(t)U_{H_{1}}(s)]}_{=0} + \underbrace{E[\alpha\beta U_{H_{1}}(t)U_{H_{2}}(s)]}_{=0} + E[\beta^{2}U_{H_{2}}(t)U_{H_{2}}(s)]$$

$$= \alpha^{2}Cov[U_{H_{1}}(t), U_{H_{1}}(s)] + \beta^{2}Cov[U_{H_{2}}(t), U_{H_{2}}(s)].$$

And the rest is clear.

(iii):

$$\begin{split} M_{\alpha,\beta,H_1,H_2}(qt) &= \alpha U_{H_1}(qt) + \beta U_{H_2}(qt) \\ &\sim \alpha q^{H_1} U_{H_1}(t) + \beta q^{H_2} U_{H_2}(t) \\ &= M_{\alpha q^{H_1},\beta q^{H_2},H_1,H_2}(t) \end{split}$$

We could give our stochastic volatility model driven by the mixed fBm. Given all parameter as in the assumption as before, we have

$$dX_{\alpha,\beta,H_1,H_2}(t) = -aX_{\alpha,\beta,H_1,H_2}(t) dt + \gamma dM_{\alpha,\beta,H_1,H_2}(t)$$
(6.31)

for $t \geq 0$ and where $a, \gamma \in \mathbb{R}_+$.

PROPOSITION 6.23. For an appropriate initial condition of (6.31), there exists a solution \hat{X}_t satisfying following properties

- (i) $(\hat{X}_t)_{t\geq 0}$ is a centered Gaussian stationary process.
- (ii) \hat{X}_t has long memory.

Proof. In terms of (6.31), then

$$X_{\alpha,\beta,H_1,H_2}(t) = X_{\alpha,\beta,H_1,H_2}(0) - a \int_0^t X_u \, du + \gamma M_{\alpha,\beta,H_1,H_2}(t).$$

Recall by (5.3), the integral in sense of Riemann-Stieljet of $\int_0^t e^{au} dM_{\alpha,\beta,H_1,H_2}$ is well-defined because $\int_0^t e^{au} dU_{H_i}$ is well-defined for $i \in \{1,2\}$.

As in Theorem 5.4, we have the solution

$$\begin{split} &X_t\\ &= e^{-at} \left(\gamma \int_0^t e^{au} \, dM_{\alpha,\beta,H_1,H_2} + X_0 \right)\\ &= e^{-at} \left(\gamma \int_0^t e^{au} \, d(\alpha U_{H_1} + \beta U_{H_2}) + X_0 \right)\\ &= e^{-at} \left(\alpha \gamma \int_0^t e^{au} \, dU_{H_1} + \beta \gamma \int_0^t e^{au} \, dU_{H_2} + X_0 \right) \end{split}$$

Given X_0 so that

$$\hat{X}_{\alpha,\beta,H_1,H_2}(t) = \alpha \underbrace{\gamma e^{-at} \int_{-\infty}^{t} e^{au} dU_{H_1}}_{:=J_{H_1}(t)} + \beta \underbrace{\gamma e^{-at} \int_{-\infty}^{t} e^{au} dU_{H_2}}_{:=J_{H_2}(t)}.$$
(6.32)

Notice $(J_{H_1}(t)), (J_{H_2}(t))$ are stationary fractional Ornstein-Uhlenbeck process. Since $J_{H_1}(t), J_{H_2}(t)$ defined as integral of U_{H_1}, U_{H_2} of Riemann-Stieljes sense, they are therefore independent. Hence

$$\sum_{k=1}^{d} c_k \hat{X}_{\alpha,\beta,H_1,H_2}$$

$$= \sum_{k=1}^{d} c_k \alpha J_{H_1}(k) + \sum_{k=1}^{d} c_k \beta J_{H_2}(k)$$

are centered Gaussian.

For (i):

$$\begin{split} & \quad \mathrm{E}[\exp(i\xi\sum_{k=1}^{d}c_{k}X_{\alpha,\beta,H_{1},H_{2}}(k))] \\ & = \quad \mathrm{E}[\exp(i\xi\sum_{k=1}^{d}c_{k}(\alpha J_{H_{1}}(k)+\beta J_{H_{2}}(k)))] \\ & = \quad \mathrm{E}[\exp(i\xi\sum_{k=1}^{d}c_{k}(\alpha J_{H_{1}}(k)))]\mathrm{E}[\exp(i\xi\sum_{k=1}^{d}c_{k}(\beta J_{H_{2}}(k)))] \\ & = \quad \mathrm{E}[\exp(i\xi\sum_{k=1}^{d}c_{k}(\alpha J_{H_{1}}(k+s)))]\mathrm{E}[\exp(i\xi\sum_{k=1}^{d}c_{k}(\beta J_{H_{2}}(k+s)))] \\ & = \quad \mathrm{E}[\exp(i\xi\sum_{k=1}^{d}c_{k}(\alpha J_{H_{1}}(k+s)+\beta J_{H_{2}}(k+s)))] \end{split}$$

for a fixed s. This shows $(\hat{X}_{\alpha,\beta,H_1,H_2}(t))$ is stationary. For (ii):

$$\varsigma(\tau)
= \operatorname{Cov}[\hat{X}_{\alpha,\beta,H_{1},H_{2}}(0), \hat{X}_{\alpha,\beta,H_{1},H_{2}}(\tau)]
= \operatorname{E}[\hat{X}_{\alpha,\beta,H_{1},H_{2}}(0)\hat{X}_{\alpha,\beta,H_{1},H_{2}}(\tau)]
= \operatorname{E}[(\alpha J_{H_{1}}(0) + \beta J_{H_{2}}(0))(\alpha J_{H_{1}}(\tau) + \beta J_{H_{2}}(\tau)]
= \alpha^{2} \operatorname{E}[J_{H_{1}}(0)J_{H_{1}}(\tau)] + \alpha\beta \operatorname{E}[J_{H_{1}}(0)J_{H_{2}}(\tau)] + \alpha\beta \operatorname{E}[J_{H_{2}}(0)J_{H_{1}}(\tau)] + \beta^{2} \operatorname{E}[J_{H_{2}}(0)J_{H_{2}}(\tau)].$$

They are centered Gaussian, using independence again, we have

 $\varsigma(\tau)$ tends to zero as τ goes to zero. As by Corollary 5.7, the last line of the equation diverges, when $H_2 \in (\frac{1}{2}, 1)$. Thus, (\hat{X}_t) has long memory property.

DEFINITION 6.24. We add a restriction for $\alpha > 0, \beta > 0$ such that $\alpha^2 + \beta^2 = 1$ to (6.31) and let a so small such that it close to zero, then we get our weighted fractional Brownian motion.

PROPOSITION 6.25. Let M_{α,β,H_1,H_2} be a weighted fractional brownian motion with respect to U_{H_1} and U_{H_2} , $T, \tau > 0$, a, γ are defined by (6.31), J_{H_1}, J_{H_2} are defined by (6.32), $\phi = H_1 - \frac{1}{2}, \psi = H_2 - \frac{1}{2}$. Then, for $t \in [0,T]$,

- (i) $E[\sup_{t \in [0,T]} |\hat{X}_{\alpha,\beta,H_1,H_2}(t) U_{H_1}(t)|] \to 0 \text{ as } a \to 0, \alpha \to 1.$
- (ii) $E[|\hat{X}_{\alpha,\beta,H_1,H_2}(t+\tau) \hat{X}_{\alpha,\beta,H_1,H_2}(t)|^2] \to \gamma^2 \tau^{2H} \text{ as } a \to 0, \alpha \to 1.$
- (iii) Let $n \in \mathbb{N}$, define

$$\begin{split} &\tilde{X}_{\alpha,\beta,H_1,H_2}(t) \\ &:= & \alpha \sum_{k=1}^{\lfloor nt \rfloor} \frac{(t - \frac{k-1}{n})^{\phi}}{\Gamma(1+\phi)} \left(J_{H_1}^{(\phi)}(\frac{k}{n}) - J_{H_1}^{(\phi)}(\frac{k-1}{n}) \right) \\ &+ & \beta \sum_{k=1}^{\lfloor nt \rfloor} \frac{(t - \frac{k-1}{n})^{\psi}}{\Gamma(1+\psi)} \left(J_{H_2}^{(\psi)}(\frac{k}{n}) - J_{H_2}^{(\psi)}(\frac{k-1}{n}) \right) \end{split}$$

then, as n goes to infinity,

$$\tilde{X}_{\alpha,\beta,H_1,H_2}(t) \to \hat{X}_{\alpha,\beta,H_1,H_2}(t)$$

in distribution.

Proof. (i):

$$\begin{split} & & \quad \mathbf{E}[\sup_{t\in[0,T]}|\hat{X}_{\alpha,\beta,H_1,H_2}(t)-U_{H_1}|]\\ = & \quad \mathbf{E}[\sup_{t\in[0,T]}|\hat{X}_{\alpha,\beta,H_1,H_2}(t)-\alpha J_{H_1}(t)+\alpha J_{H_1}(t)-U_{H_1}(t)|]\\ \leq & \quad \mathbf{E}[\sup_{t\in[0,T]}|\underbrace{\hat{X}_{\alpha,\beta,H_1,H_2}(t)-\alpha J_{H_1}(t)|}_{\overset{\alpha\uparrow^1}{\to}0}] + \mathbf{E}[\sup_{t\in[0,T]}|\underbrace{\alpha J_{H_1}(t)}_{\overset{\alpha\uparrow^1}{\to}J_{H_1}(t)}-U_{H_1}(t)|]\\ \overset{\alpha\uparrow^1}{\to} & \quad \mathbf{E}[\sup_{t\in[0,T]}|J_{H_1}(t)-U_{H_1}(t)|]\\ \overset{\alpha\downarrow 0}{\to} & \quad 0 \end{split}$$

(ii):

$$E[|\hat{X}_{\alpha,\beta,H_{1},H_{2}}(t+\tau) - \hat{X}_{\alpha,\beta,H_{1},H_{2}}(t)|^{2}]$$

$$= 2Var[\hat{X}_{\alpha,\beta,H_{1},H_{2}}(t)] - 2Cov[\hat{X}_{\alpha,\beta,H_{1},H_{2}}(t), M_{\alpha,\beta,H_{1},H_{2}}(t+\tau)]$$

$$= 2(\alpha^{2}Var[J_{H_{1}}(t)] + \beta^{2}Var[J_{H_{2}}(t)]$$

$$- \alpha^{2}Cov[J_{H_{1}}(t), J_{H_{1}}(t+\tau)] - \beta^{2}Cov[J_{H_{2}}(t), J_{H_{2}}(t+\tau)])$$

$$= 2(\alpha^{2}Var[J_{H_{1}}(t)] - \alpha^{2}Cov[J_{H_{1}}(t), J_{H_{1}}(t+\tau)]$$

$$+ \beta^{2}Var[J_{H_{2}}(t)] - \beta^{2}Cov[J_{H_{2}}(t), J_{H_{2}}(t+\tau)])$$

$$\rightarrow 2(\alpha^{2}(\frac{1}{2}\gamma^{2}\tau^{2H_{1}}) + \beta^{2}\varsigma_{\hat{J}_{H_{2}}}(\tau))$$

$$\stackrel{\alpha\uparrow 1}{\rightarrow} \gamma^{2}\tau^{2H_{1}}$$

(iii): Suppose $\{Y_t\}_t$, $\{Z_t\}_t$ are two families of random variables with $Y_t \to Y, Z_t \to Z$ in distribution. Y_t, Z_t are independent for each t. Then $Y_t + Z_t \to Y + Z$, because, using continuity theorem of characteristic function

$$E[\exp i\xi(Y_t + Z_t)] = E[\exp i\xi Y_t]E[\exp i\xi Z_t]$$

$$\to E[\exp i\xi Y]E[\exp i\xi Z]$$

$$= E[\exp i\xi(Y + Z)].$$

Consider

$$\alpha \sum_{k=1}^{\lfloor nt \rfloor} \frac{(t-\frac{k-1}{n})^{\phi}}{\Gamma(1+\phi)} \left(J_{H_1}^{(\phi)}(\frac{k}{n}) - J_{H_1}^{(\phi)}(\frac{k-1}{n}) \right)$$

and

$$\beta \sum_{k=1}^{\lfloor nt \rfloor} \frac{(t-\frac{k-1}{n})^{\psi}}{\Gamma(1+\psi)} \left(J_{H_2}^{(\psi)}(\frac{k}{n}) - J_{H_2}^{(\psi)}(\frac{k-1}{n}) \right)$$

are independent. According Theorem 6.17 and Theorem of Cramer Slutsky.

$$\alpha \sum_{k=1}^{\lfloor nt \rfloor} \frac{(t - \frac{k-1}{n})^{\phi}}{\Gamma(1+\phi)} \left(J_{H_1}^{(\phi)}(\frac{k}{n}) - J_{H_1}^{(\phi)}(\frac{k-1}{n}) \right) \to \alpha J_{H_1}(t)$$

and

$$\beta \sum_{k=1}^{\lfloor nt \rfloor} \frac{(t - \frac{k-1}{n})^{\psi}}{\Gamma(1 + \psi)} \left(J_{H_2}^{(\psi)}(\frac{k}{n}) - J_{H_2}^{(\psi)}(\frac{k-1}{n}) \right) \to \beta J_{H_2}(t)$$

in distribution. Then, since mentioned above,

$$\tilde{X}_{\alpha,\beta,H_1,H_2}(t) \rightarrow \alpha J_{H_1}(t) + \beta J_{H_2}(t)
= \hat{X}_{\alpha,\beta,H_1,H_2}(t)$$

6.6 Discussion

In order to model log-volatility, we take fOU process into account. In FSV, we choose $H > \frac{1}{2}$ and that will make sure the solution of log-volatility SDE has long memory. In contrast, although it could not exhibit the long memory by $H < \frac{1}{2}$, RFSV demonstrates a more reasonable smoothness of volatility. For instance, RFSV ensures the slop of the plotting of $log(s(\tau, \hat{X}))$ against $log(\tau, \hat{X})$, which consists with the empirical result we observed, see (6.27).

The weighted-FSV model inherits long memory of FSV. With an adjustable factor, one can achieve a result of smoothness of volatility close to it by RFSV. As for RFSV, when a goes to zero, the \hat{X} acts locally as fBm at any compact time scale.

Not only in FSV but also in RFSV, there is a discretization with the fractional derivative, which is an AR(1) process. However, it is not the case by weighted-FSV, in which \hat{X} is approached as the sum of two fractional derivatives.

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ERKLÄRUNG

Hiermit erkläre ich, dass ich die am heutigen Tag eingereichte Diplomarbeit zum Thema "Fractional Brownian motion and applications in financial mathematics" unter Betreuung von Prof. Dr. rer. nat. M. Keller-Ressel selbstständig erarbeitet, verfasst und Zitate kenntlich gemacht habe. Andere als die angegebenen Hilfsmittel wurde von mir nicht benutzt.

Datum Unterschrift