

Path Independence of Additive Functionals for SDEs under G -framework

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Abstract

The path independence of additive functionals for SDEs driven by the G -Brownian motion is characterized by nonlinear PDEs. The main result generalizes the existing ones for SDEs driven by the standard Brownian motion.

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1 Introduction

Stochastic differential equations (SDEs) under the linear probability space have been widely used in modeling financial markets and economic phenomena [1, 2]. However, in many practical situations, most of the financial activities take place with uncertainty [3], for which a fundamental theory of SDEs driven by the G -Brownian motion (G -SDEs) has been developed in [11, 12, 13]. Since then G -SDEs have received much attention, see for instance [8] on the Feynman-Kac formula, [9, 10] on the stochastic control, [5, 6] on the ergodicity, [21, 25] on the stochastic stability, and [7] on the G -SPDEs.

In the equilibrium financial market, there exists a risk neutral measure which admits a path independent density with respect to the real world probability [27]. To construct such risk neutral measures, the path independence of additive functionals for SDEs has been investigated extensively; see [23] for the pioneer work. Subsequently, [23] has been

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extended in [16, 17, 26] for finite dimensional SDEs, and in [18, 24] for infinite dimensional SPDEs, where [26] allows the SDEs involved to be degenerate. Recently, [19] investigated the path independence of additive functionals for a class of distribution dependent SDEs. Nevertheless, all of these papers only focus on linear probability spaces. To fill this gap, in this paper, we intend to characterize the path independence of additive functionals for G -SDEs. To this end, below we recall some basic facts on the G -Brownian motion.

For a positive integer d , let $(\mathbb{R}^d, \langle \cdot, \cdot \rangle, |\cdot|)$ be the d -dimensional Euclidean space, $\mathbb{R}^d \otimes \mathbb{R}^d$ the family of all $d \times d$ -matrices, \mathbb{S}^d the collection of all symmetric $d \times d$ -matrices, $\mathbf{0}_d \in \mathbb{R}^d$ the zero vector, $\mathbf{0}_{d \times d} \in \mathbb{R}^d \otimes \mathbb{R}^d$ the zero matrix, and $\mathbf{I}_{d \times d} \in \mathbb{R}^d \otimes \mathbb{R}^d$ the identity matrix. For a matrix A , let A^* be its transpose and $\|A\|_{\text{HS}} = (\text{trace}(AA^*))^{1/2}$ be its Hilbert-Schmidt norm (or Frobenius norm). For a number $a \in \mathbb{R}$, a^+ and a^- stipulate its positive part and negative part, respectively. For $\sigma_1, \sigma_2 \in \mathbb{S}^d$, the notation $\sigma_1 \leq \sigma_2$ (res. $\sigma_1 < \sigma_2$) means that $\sigma_2 - \sigma_1$ is non-negative (res. positive) definite, and we let

$$[\sigma_1, \sigma_2] := \{\gamma | \gamma \in \mathbb{S}^d, \sigma_1 \leq \gamma \leq \sigma_2\}.$$

Let $C^{1,2}(\mathbb{R}_+ \times \mathbb{R}^d; \mathbb{R})$ be the collection of all continuous functions $V : \mathbb{R}_+ \times \mathbb{R}^d \rightarrow \mathbb{R}$ which are once differentiable w.r.t. the first argument, twice differentiable w.r.t. the second argument, and all these derivatives are joint continuous. Write ∇ and ∇^2 by the gradient operator and Hessian operator, respectively.

For any fixed $T > 0$,

$$\Omega_T = \{\omega | [0, T] \ni t \mapsto \omega_t \in \mathbb{R}^d \text{ is continuous with } \omega(0) = \mathbf{0}_d\}$$

endowed with the uniform topology. Let $B_t(\omega) = \omega_t, \omega \in \Omega_T$, be the canonical process. Set

$$L_{ip}(\Omega_T) := \{\varphi(B_{t_1}, \dots, B_{t_n}), n \in \mathbb{N}, t_1, \dots, t_n \in [0, T], \varphi \in C_{b, lip}(\mathbb{R}^d \otimes \mathbb{R}^n)\},$$

where $C_{b, lip}(\mathbb{R}^d \otimes \mathbb{R}^n)$ denotes the set of bounded Lipschitz functions $f : \mathbb{R}^d \otimes \mathbb{R}^n \rightarrow \mathbb{R}$. Let $G : \mathbb{S}^d \rightarrow \mathbb{R}$ be a monotonic, sublinear and homogeneous function; see e.g. [13, p16]. Throughout the paper, we always assume that $G : \mathbb{S}^d \rightarrow \mathbb{R}$ is non-degenerate, i.e., there exists some $\delta > 0$ such that

$$(1.1) \quad G(A) - G(B) \geq \frac{\delta}{2} \text{trace}[A - B], \quad A \geq B, A, B \in \mathbb{S}^d.$$

For any $\xi \in L_{ip}(\Omega_T)$, i.e.,

$$\xi(\omega) = \varphi(\omega(t_1), \dots, \omega(t_n)), \quad 0 = t_0 < t_1 < \dots < t_n = T,$$

the conditional G -expectation is defined by

$$\bar{\mathbb{E}}_t[\xi] := u_k(t, \omega(t); \omega(t_1), \dots, \omega(t_{k-1})), \quad \xi \in L_{ip}(\Omega_T), \quad t \in [t_{k-1}, t_k], \quad k = 1, \dots, n,$$

where $(t, x) \mapsto u_k(t, x; x_1, \dots, x_{k-1})$, $k = 1, \dots, n$, solves the following G -heat equation

$$(1.2) \quad \begin{cases} \partial_t u_k + G(\partial_x^2 u_k) = 0, & (t, x) \in [t_{k-1}, t_k] \times \mathbb{R}^d, \quad k = 1, \dots, n, \\ u_k(t_k, x; x_1, \dots, x_{k-1}) = u_{k+1}(t_k, x; x_1, \dots, x_{k-1}, x_k), & k = 1, \dots, n-1, \\ u_n(t_n, x; x_1, \dots, x_{n-1}) = \varphi(x_1, \dots, x_{n-1}, x), & k = n. \end{cases}$$

Since G is non-degenerate, the solution of (1.2) satisfies $u_k \in C^{1,2}(\mathbb{R}_+ \times \mathbb{R}^d; \mathbb{R})$; see [13, Appendix C, Theorem 4.5, p127]. The corresponding G -expectation of ξ is defined by $\bar{\mathbb{E}}[\xi] = \bar{\mathbb{E}}_0[\xi]$. Then the canonical process $B_t(\omega) := \omega_t$ is called a G -Brownian motion in $(\Omega_T, L_G^p(\Omega_T), \bar{\mathbb{E}})$, where $L_G^p(\Omega_T)$ is the completion of $L_{ip}(\Omega_T)$ under the norm $(\bar{\mathbb{E}}[|\cdot|^p])^{\frac{1}{p}}$, $p \geq 1$. By definition, we have $G(A) = \frac{1}{2}\bar{\mathbb{E}}\langle AB_1, B_1 \rangle$, $A \in \mathbb{S}^d$. The function G is called the generating function corresponding of the d -dimensional G -Brownian motion $(B_t)_{t \geq 0}$. According to [13], there exists a bounded, convex, and closed subset $\Theta \subset \mathbb{S}^m$ such that

$$(1.3) \quad G(A) = \frac{1}{2} \sup_{Q \in \Theta} \text{trace}[AQ], \quad A \in \mathbb{S}^d.$$

In particular, for 1-dimensional G -Brownian motion $(B_t)_{t \geq 0}$, one has $G(a) = (\bar{\sigma}^2 a^+ - \underline{\sigma}^2 a^-)/2$, $a \in \mathbb{R}$, where $\bar{\sigma}^2 := \bar{\mathbb{E}}[B_1^2] \geq -\bar{\mathbb{E}}[-B_1^2] =: \underline{\sigma}^2 > 0$.

Let

$$M_G^{p,0}([0, T]) = \left\{ \eta_t := \sum_{j=0}^{N-1} \xi_j I_{[t_j, t_{j+1})}(t); \xi_j \in L_G^p(\Omega_{t_j}), N \in \mathbb{N}, 0 = t_0 < t_1 < \dots < t_N = T \right\}.$$

Let $M_G^p([0, T])$ and $H_G^p([0, T])$ be the completion of $M_G^{p,0}([0, T])$ under the norm

$$\|\eta\|_{M_G^p([0, T])} := \left(\bar{\mathbb{E}} \int_0^T |\eta_t|^p dt \right)^{\frac{1}{p}}, \quad \|\eta\|_{H_G^p([0, T])} := \left\{ \bar{\mathbb{E}} \left(\int_0^T |\eta_t|^2 dt \right)^{\frac{p}{2}} \right\}^{\frac{1}{p}},$$

respectively. We need to point out that if $p = 2$, then $M_G^p([0, T]) = H_G^p([0, T])$. Denote by $M_G^p([0, T]; \mathbb{R}^d)$ all d -dimensional stochastic processes $\eta_t = (\eta_t^1, \dots, \eta_t^d)$, $t \geq 0$ with $\eta_t^i \in M_G^p([0, T])$. Let $H_G^1([0, T]; \mathbb{R}^d)$ be all d -dimensional stochastic processes $\zeta_t = (\zeta_t^1, \dots, \zeta_t^d)$, $t \geq 0$ with $\zeta^i \in H_G^1([0, T])$.

Furthermore, we also need the Choquet capacity associated with the G -expectation. Let \mathcal{M} be the collection of all probability measures on $(\Omega_T, \mathcal{B}(\Omega_T))$. According to [4], there exists a weakly compact subset $\mathcal{P} \subset \mathcal{M}$ such that

$$\bar{\mathbb{E}}[X] = \sup_{P \in \mathcal{P}} \mathbb{E}_P[X], \quad X \in L_{ip}(\Omega_T).$$

Then the associated Choquet capacity is defined by

$$c(A) = \sup_{P \in \mathcal{P}} P(A), \quad A \in \mathcal{B}(\Omega_T).$$

A set $A \subset \Omega_T$ is called polar if $c(A) = 0$, and we say that a property holds quasi-surely (q.s.) if it holds outside a polar set.

In this paper, we consider the following G -SDE

$$(1.4) \quad dX_t = b(t, X_t)dt + \sum_{i,j=1}^d h_{ij}(t, X_t) d\langle B^i, B^j \rangle_t + \langle \sigma(t, X_t), dB_t \rangle,$$

where $b, h_{ij} = h_{ji} : [0, T] \times \mathbb{R}^d \rightarrow \mathbb{R}^d$ and $\sigma : [0, T] \times \mathbb{R}^d \rightarrow \mathbb{R}^d \otimes \mathbb{R}^d$, B_t is a d -dimensional G -Brownian motion, and $\langle B^i, B^j \rangle_t$ stands for the mutual variation process of the i -th component B_t^i and the j -th component B_t^j . To ensure the existence and uniqueness of the solution of (1.4) in $M_G^2([0, T]; \mathbb{R}^d)$, we assume

$$(1.5) \quad |b(t, x) - b(t, y)| + \sum_{i,j=1}^d |h_{ij}(t, x) - h_{ij}(t, y)| + \|\sigma(t, x) - \sigma(t, y)\|_{\text{HS}} \leq K|x - y|,$$

for some constant $K \geq 0$ and all $t \in [0, T]$, $x, y \in \mathbb{R}^d$; see [13, Theorem 1.2, p82].

Now we recall from [19] the following notions for the path independence of additive functionals.

Definition 1.1. For $f = (f_{ij}) : \mathbb{R}_+ \times \mathbb{R}^d \rightarrow \mathbb{S}^d$ and $g : \mathbb{R}_+ \times \mathbb{R}^d \rightarrow \mathbb{R}^d$, the additive functional $(A_{s,t}^{f,g})_{0 \leq s \leq t}$ is defined by

$$(1.6) \quad A_{s,t}^{f,g} = \beta \int_s^t G(f)(r, X_r) dr + \alpha \sum_{i,j=1}^d \int_s^t f_{ij}(r, X_r) d\langle B^i, B^j \rangle_r + \int_s^t \langle g(r, X_r), dB_r \rangle,$$

where $\beta, \alpha \in \mathbb{R}$ are two parameters, $f = f^*$, and $(X_r)_{r \geq 0}$ solves (1.4).

Definition 1.2. The additive functional $A_{s,t}^{f,g}$ is said to be path independent, if there exists a function $V : \mathbb{R}_+ \times \mathbb{R}^d \rightarrow \mathbb{R}$ such that for any $s \in [0, T]$ and any solution $(X_t)_{t \in [s, T]}$ to (1.4) from time s , it holds

$$(1.7) \quad A_{s,t}^{f,g} = V(t, X_t) - V(s, X_s), \quad t \in [s, T].$$

In terms of **Definition 1.2**, the path independence of the additive functional $A_{s,t}^{f,g}$ means that $A_{s,t}^{f,g}$ depends only on X_s and X_t but not the path $(X_r)_{s < r < t}$, for any solution $(X_r)_{r \in [s, T]}$ to (1.4) from times s and any $t \in (s, T]$.

The aim of this paper is to provide sufficient and necessary characterizations for the path independence of the additive functional $A_{s,t}^{f,g}$.

To see that (1.6) covers additive functionals investigated in existing references for the path independence under the linear probability space, let Θ in (1.3) be a singleton: $\Theta = Q$, and $\bar{\mathbb{E}} = \mathbb{E}$ be a linear expectation. Then the associated G -Brownian motion B_t becomes the classical zero-mean normal distributed with covariance Q . Specially, let $\sigma^2 = \bar{\sigma}^2 = \mathbf{1}_{d \times d}$, i.e., $Q = \mathbf{1}_{d \times d}$, $G(A) = \frac{1}{2} \text{trace}(A)$, $A \in \mathbb{S}^d$, we have $d\langle B^i, B^j \rangle_r = \delta_{ij} dr$, where δ_{ij} is a indicative function, $1 \leq i, j \leq d$, and (1.6) reduces to

$$A_{s,t}^{f,g} = \left(\alpha + \frac{\beta}{2} \right) \sum_{i=1}^d \int_s^t f_{ii}(r, X_r) dr + \int_s^t \langle g(r, X_r), dB_r \rangle.$$

Taking $\mathbf{f}(r, X_r) = \left(\alpha + \frac{\beta}{2} \right) \sum_{i=1}^d f_{ii}(r, X_r)$, this goes back to the additive functional studied in [19]:

$$A_{s,t}^{f,g} := \int_s^t \mathbf{f}(r, X_r) dr + \int_s^t \langle g(r, X_r), dB_r \rangle.$$

In particular, when $\mathbf{f} = \frac{1}{2}|g|^2$, we have

$$(1.8) \quad A_{s,t}^{\frac{1}{2}|g|^2,g} = \frac{1}{2} \int_s^t |g|^2(r, X_r) dr + \int_s^t \langle g(r, X_r), dB_r \rangle, \quad 0 \leq s \leq t,$$

which corresponds to the Girsanov transform $d\mathbb{Q}_{s,t} := \exp\{-A_{s,t}^{\frac{1}{2}|g|^2,g}\}d\mathbb{P}$. To make the solution X_t of (1.4) a martingale under $\mathbb{Q}_{s,t}$, we reformulate (1.4) as

$$dX_t = \{b + h\}(t, X_t)dt + \langle \sigma(t, X_t), dB_t \rangle,$$

where

$$(1.9) \quad h(t, u) := \sum_{i=1}^d h_{ii}(t, u), \quad (t, u) \in \mathbb{R}_+ \times \mathbb{R}^d.$$

When σ is invertible, taking $g = \sigma^{-1}(b + h)$ in (1.8), we have

$$A_{s,t}^{\frac{1}{2}|g|^2,g} = \frac{1}{2} \int_s^t |\sigma^{-1}(b + h)|^2(r, X_r) dr + \int_s^t \langle (\sigma^{-1}(b + h))(r, X_r), dB_r \rangle, \quad 0 \leq s \leq t.$$

Then, by the Girsanov theorem, $(X_r)_{s \leq r \leq t}$ is a martingale under $\mathbb{Q}_{s,t}$, which fits well the requirement of risk neutral measure in finance. The path independence of this particular additive functional has been investigated in [16, 19, 23, 24, 26].

Remark 1.1. When $\alpha \neq 0$, (1.6) is equivalent to

$$\begin{aligned} \alpha^{-1} A_{s,t}^{f,g} &= \alpha^{-1} \beta \int_s^t G(f)(r, X_r) dr + \sum_{i,j=1}^m \int_s^t f_{ij}(r, X_r) d\langle B^i, B^j \rangle_r \\ &\quad + \int_s^t \langle \alpha^{-1} g(r, X_r), dB_r \rangle. \end{aligned}$$

So, in this case, the path independence of the additive functional (1.6) can be reduced to the case of $\alpha = 1$. However, the case for $\alpha = 0$ also includes interesting examples (see Example 4.1 below), so it is reasonable to consider $A_{s,t}^{f,g}$ in (1.6) with two parameters α and β .

The remainder of the paper is organized as follows. In Section 2, following the line of [15, 22], we present a decomposition theorem for multidimensional G -semimartingales. In Section 3, we characterize the path independence of $A_{s,t}^{f,g}$ using nonlinear PDEs, so that main results in [16, 23, 24, 26] are extended to the present nonlinear expectation setting. Finally, in Section 4, we provide an example to illustrate the main result for $\alpha = 0$ as mentioned in Remark 1.1.

2 A Decomposition Theorem

This part is essentially due to [15, 22]. Set $\delta_n(t) := \sum_{i=1}^{n-1} (-1)^i \mathbf{1}_{(\frac{i}{n}, \frac{i+1}{n}]}(t)$, $t \in [0, T]$. For any $A, B \in \mathbb{S}^d$, let $(A, B)_{\text{HS}} = \text{trace}(AB)$ and $\|A\|_{\text{HS}} = \sqrt{(A, A)_{\text{HS}}}$. Then $(\mathbb{S}^d, \langle \cdot, \cdot \rangle_{\text{HS}}, \|\cdot\|_{\text{HS}})$ is a Hilbert space; see e.g. [20]. Let $\text{spec}(\cdot)$ be the spectrum of a matrix \cdot , and let $\langle B \rangle_t = (\langle B^i, B^j \rangle_t)_{ij}$.

From now on, we consider

$$(2.1) \quad G(A) := \frac{1}{2} \sup_{\gamma \in [\underline{\sigma}, \bar{\sigma}]} \langle \gamma^2, A \rangle_{\text{HS}}, \quad A \in \mathbb{S}^d, 0 < \underline{\sigma} < \bar{\sigma} \text{ are two matrices in } \mathbb{S}^d.$$

Consequently, $\underline{\sigma}^2 < \frac{d}{dt} \langle B \rangle_t \leq \bar{\sigma}^2$, and (1.1) holds for $\delta = \lambda_0(\underline{\sigma}^2)$, where $\lambda_0(\underline{\sigma}^2) = \min\{\lambda \in \text{spec}(\underline{\sigma}^2)\}$.

Let $M_G^1([0, T]; \mathbb{S}^d)$ be all symmetric $d \times d$ matrices $\eta_t = (\eta_t^{ij})_{d \times d}$ with $\eta^{ij} \in M_G^1([0, T])$, and $\|\eta\|_{M_G^1([0, T]; \mathbb{S}^d)} = \bar{\mathbb{E}} \int_0^T \|\eta_t\|_{\text{HS}} dt$.

Let $c_0 = \min\{\lambda \in \text{spec}((\bar{\sigma}^2 - \underline{\sigma}^2)/2)\}$, $C_0 = \frac{1}{2} \|\bar{\sigma}^2 - \underline{\sigma}^2\|_{\text{HS}}$.

To make the content self-contained, we cite from [15] some well-known results and restated them as follows.

Lemma 2.1. Let G be in (2.1). For any $\eta \in M_G^1([0, T]; \mathbb{S}^d)$, the limit

$$\|\eta\|_{\mathbb{M}_G} := \lim_{n \rightarrow \infty} \bar{\mathbb{E}} \int_0^T \delta_n(s) (\eta_s, d\langle B \rangle_s)_{\text{HS}}$$

exists. Moreover, $\|\cdot\|_{\mathbb{M}_G}$ defines a norm on $M_G^1([0, T]; \mathbb{S}^d)$, and for any $0 < \epsilon \leq c_0$, it holds that,

$$\epsilon \|\eta\|_{M_{G_\epsilon}^1([0, T]; \mathbb{S}^d)} \leq \|\eta\|_{\mathbb{M}_G} \leq C_0 \|\eta\|_{M_G^1([0, T]; \mathbb{S}^d)},$$

where $G_\epsilon(A) := \frac{1}{2} \sup_{\gamma \in [\underline{\sigma}_\epsilon, \bar{\sigma}_\epsilon]} \langle \gamma^2, A \rangle_{\text{HS}}$, $\underline{\sigma}_\epsilon^2 := \underline{\sigma}^2 + \epsilon \mathbf{I}_{d \times d}$, and $\bar{\sigma}_\epsilon^2 := \bar{\sigma}^2 - \epsilon \mathbf{I}_{d \times d}$.

With Lemma 2.1 in hand, we have the following corollary which will play a crucial role in the analysis below.

Corollary 2.2. Let G be in (2.1), and let $\eta \in M_G^1([0, T]; \mathbb{S}^d)$, $\zeta \in M_G^1([0, T])$. If

$$\int_0^t (\eta_s, d\langle B \rangle_s)_{\text{HS}} = \int_0^t \zeta_s ds, \quad t \in [0, T],$$

then

$$\bar{\mathbb{E}} \int_0^T \|\eta_s\|_{\text{HS}} ds = \bar{\mathbb{E}} \int_0^T |\zeta_s| ds = 0.$$

Proof. According to Lemma 2.1 and [22, Theorem 3.3 (i)], we deduce that

$$\|\eta\|_{\mathbb{M}_G} = \lim_{n \rightarrow \infty} \bar{\mathbb{E}} \int_0^T \delta_n(s) (\eta_s, d\langle B \rangle_s)_{\text{HS}} = \lim_{n \rightarrow \infty} \bar{\mathbb{E}} \int_0^T \delta_n(s) \zeta_s ds = 0.$$

Recall from Lemma 2.1 that $\|\eta\|_{\mathbb{M}_G}$ is a norm, then, c -q.s., $\eta_t \equiv \mathbf{0}_{d \times d}$, a.e. $t \in [0, T]$. Therefore, $\bar{\mathbb{E}} \int_0^T \|\eta_s\|_{\text{HS}} ds = 0$, which leads to $\bar{\mathbb{E}} \int_0^T |\zeta_s| ds = 0$. \square

Consider the following Itô process in $(\Omega_T, L_G^p(\Omega_T), \bar{\mathbb{E}})$

$$(2.2) \quad X_t = \int_0^t \xi_r dr + \sum_{i,j=1}^d \int_0^t \eta_r^{ij} d\langle B^i, B^j \rangle_r + \sum_{i=1}^d \int_0^t \zeta_r^i dB_r^i, \quad t \geq 0,$$

where $\xi, \eta^{ij} \in M_G^1([0, T]; \mathbb{R}^d)$ with $\eta^{ij} = \eta^{ji}$, and $\zeta^i \in H_G^1([0, T]; \mathbb{R}^d)$.

Now we can state the following decomposition theorem.

Theorem 2.3. *For G in (2.1) and let X_t be in (2.2). Then $X_t = \mathbf{0}_d$ for all $t \in [0, T]$ if and only if on $\Omega_T \times [0, T]$ it holds $c \times dt - q.s. \times a.e.$, $\xi_t = \mathbf{0}_d, \eta_t^{ij} = \mathbf{0}_d, \zeta_t^i = \mathbf{0}_d, i, j = 1, \dots, d$.*

Proof. The proof of the sufficiency is trivial, it suffices to prove the necessity. Assume $X_t = \mathbf{0}_d$ for $t \in [0, T]$. Then (2.2) is equivalent to

$$(2.3) \quad \int_0^t \xi_s^k ds + \sum_{i,j=1}^d \int_0^t \eta_r^{kij} d\langle B^i, B^j \rangle_r + \sum_{j=1}^d \int_0^t \zeta_s^{kj} dB_s^j = \mathbf{0}_d, \quad k = 1, \dots, d, \quad t \in [0, T],$$

where ξ^k (resp. η^{kij}) denotes the k -th component of the column vector ξ (resp. η^{ij}). Taking quadratic processes w.r.t. $\int_0^\cdot \zeta_s^{ki} dB_s^i$ on both side of (2.3), we deduce that

$$\begin{aligned} 0 &= \sum_{i,j=1}^d \left\langle \int_0^\cdot \zeta_s^{kj} dB_s^j, \int_0^\cdot \zeta_s^{ki} dB_s^i \right\rangle_t = \left\langle \sum_{i=1}^d \int_0^\cdot \zeta_s^{ki} dB_s^i \right\rangle_t \\ &= \sum_{i,j=1}^d \int_0^t \zeta_s^{kj} \zeta_s^{ki} d\langle B^i, B^j \rangle_s = \int_0^t ((\zeta_s^k)^* \zeta_s^k, d\langle B \rangle_s)_{\text{HS}} \\ &= \int_0^t \langle d\langle B \rangle_s (\zeta_s^k)^*, (\zeta_s^k)^* \rangle \geq \int_0^t \langle \sigma^2 (\zeta_s^k)^*, (\zeta_s^k)^* \rangle ds \geq 0 \end{aligned}$$

with $\zeta^k := (\zeta^{k1}, \dots, \zeta^{kd})$. Since $\sigma^2 > 0$, this implies $\zeta_t = \mathbf{0}_{d \times d}$, a.e. $t \in [0, T]$.

It remains to show that $\xi_t = \mathbf{0}_d$ and $\eta_t = \mathbf{0}_d$. In fact, since $\zeta_t = \mathbf{0}_d$, we have

$$\int_0^t -\xi_s^k ds = \sum_{i,j=1}^d \int_0^t \eta_s^{kij} d\langle B^i, B^j \rangle_s = \int_0^t (\eta_s^k, d\langle B \rangle_s)_{\text{HS}}, \quad k = 1, \dots, d, \quad t \in [0, T]$$

with $\eta^k = (\eta^{kij})_{ij}$. By Corollary 2.2, this implies

$$\bar{\mathbb{E}} \int_0^T |\xi_s^k| ds = \bar{\mathbb{E}} \int_0^T \|\eta_s^k\|_{\text{HS}} ds = 0, \quad k = 1, \dots, d.$$

Thus, we conclude that c -q.s. for a.e. $t \in [0, T]$, $\eta_t^k = \mathbf{0}_{d \times d}$ and $\xi_t^k = 0$, $k = 1, \dots, d$. Therefore, $\xi_t = \eta_t^{ij} = \mathbf{0}_d$.

□

3 Characterization of Path Independence

The main result of the paper is the following.

Theorem 3.1. *Let G be in (2.1). Then $A_{s,t}^{f,g}$ is path independent in the sense of (1.7) for some $V \in C^{1,2}(\mathbb{R}_+ \times \mathbb{R}^d)$ if and only if*

$$(3.1) \quad \begin{cases} \frac{\partial}{\partial t}V(t, x) = \beta G(f)(t, x) - \langle \nabla V, b \rangle(t, x), \\ \alpha f_{ij}(t, x) = \left(\langle \nabla V, h_{ij} \rangle + \frac{1}{2} \langle \sigma_i, (\nabla^2 V) \sigma_j \rangle \right)(t, x), \\ g(t, x) = (\sigma^* \nabla V)(t, x), \quad (t, x) \in [0, T] \times \mathbb{R}^d, \end{cases}$$

where $i, j = 1, \dots, d$, and σ_i stands for the i -th column of σ .

Proof. We first prove the necessity. For any $(s, x) \in [0, T] \times \mathbb{R}^d$, let $(X_t)_{t \geq s}$ solves (1.4) with $X_s = x$. Since $(A_{s,t}^{f,g})_{t \in [s, T]}$ is path independent in the sense of (1.7), it follows that

$$(3.2) \quad \begin{aligned} dV(t, X_t) &= \beta G(f)(t, X_t)dt + \alpha \sum_{i,j=1}^d f_{ij}(t, X_t) d\langle B^i, B^j \rangle_t \\ &\quad + \langle g(t, X_t), dB_t \rangle, \quad t \in [s, T]. \end{aligned}$$

On the other hand, by Itô's formula, we derive that

$$(3.3) \quad \begin{aligned} dV(t, X_t) &= \left(\frac{\partial}{\partial t}V + \langle \nabla V, b \rangle \right)(t, X_t)dt + \langle (\sigma^* \nabla V)(t, X_t), dB_t \rangle \\ &\quad + \sum_{i,j=1}^d \left(\langle \nabla V, h_{ij} \rangle + \frac{1}{2} \langle \sigma_i, (\nabla^2 V) \sigma_j \rangle \right)(t, X_t) d\langle B^i, B^j \rangle_t, \quad t \in [s, T]. \end{aligned}$$

Since coefficients b, h and σ satisfy the Lipschitz condition in (1.5), and the solution of (1.4) satisfies $X_t \in M_G^2([0, T]; \mathbb{R}^d)$, it's not difficult to verify $\left(\frac{\partial}{\partial t}V + \langle \nabla V, b \rangle \right)(t, X_t) \in M_G^1([0, T])$, $\left(\langle \nabla V, h_{ij} \rangle + \frac{1}{2} \langle \sigma_i, (\nabla^2 V) \sigma_j \rangle \right)(t, X_t) \in M_G^1([0, T])$, and $(\sigma_i^* \nabla V)(t, X_t) \in H_G^1([0, T])$, thus hypotheses of Theorem 2.3 are satisfied. Combining (3.2) and (3.3), and applying Theorem 2.3 for the process $V(t, X_t)$, we obtain c -q.s. for a.e. $t \in [s, T]$,

$$(3.4) \quad \begin{cases} \left(\frac{\partial}{\partial t}V + \langle \nabla V, b \rangle \right)(t, X_t) = \beta G(f)(t, X_t), \\ \alpha f_{ij}(t, X_t) = \left(\langle \nabla V, h_{ij} \rangle + \frac{1}{2} \langle \sigma_i, (\nabla^2 V) \sigma_j \rangle \right)(t, X_t), \\ g(t, X_t) = (\sigma^* \nabla V)(t, X_t). \end{cases}$$

Since all terms in (3.4) are continuous in t , these equations hold c -q.s. at $t = s$, so by $X_s = x$, we have

$$(3.5) \quad \begin{cases} \left(\frac{\partial}{\partial t}V + \langle \nabla V, b \rangle \right)(s, x) = \beta G(f)(s, x), \\ \alpha f_{ij}(s, x) = \left(\langle \nabla V, h_{ij} \rangle + \frac{1}{2} \langle \sigma_i, (\nabla^2 V) \sigma_j \rangle \right)(s, x), \\ g(s, x) = (\sigma^* \nabla V)(s, x). \end{cases}$$

Due to the arbitrariness of s and x , we prove (3.1).

Next, for the sufficiency, taking advantage of (3.1), we deduce from (3.3) that (3.2) holds true. By taking stochastic integration we prove (1.7), and therefore complete the proof. \square

Let us comparison this result with known ones in the linear expectation setting.

Remark 3.2. *Comparing with the Girsanov transform in the linear expectation setting as mentioned in Introduction, we take for instance $\alpha = 1, \beta = -1, h_{ij} = 0$ and*

$$(3.6) \quad f_{ii} = \frac{1}{d}|g|^2, \quad 1 \leq i \leq d.$$

When $\underline{\sigma}^2 = \bar{\sigma}^2 = \mathbf{1}_{d \times d}$, this goes back to the classic linear expectation, $(B_t)_{t \geq 0}$ is a d -dimensional standard Brownian motion defined on the probability space $(\Omega, \mathcal{F}, \mathbb{P})$, we have $\langle B^i, B^j \rangle_r = \delta_{ij}r$, and

$$G(f) = \frac{|g|^2}{2d} \text{trace}[\mathbf{1}_{d \times d}] = \frac{|g|^2}{2}.$$

So

$$\begin{aligned} A_{s,t}^{f,g} : &= \beta \int_s^t G(f)(r, X_r) dr + \alpha \sum_{i,j=1}^d \int_s^t f_{ij}(r, X_r) d\langle B^i, B^j \rangle_r \\ &\quad + \int_s^t \langle g(r, X_r), dB_r \rangle \\ &= \frac{1}{2} \int_s^t |g(r, X_r)|^2 dr + \int_s^t \langle g(r, X_r), dB_r \rangle \end{aligned}$$

gives the weighted probability $\exp\{-A_{s,t}^{f,g}\}d\mathbb{P}$ in the Girsanov theorem.

By taking $\alpha = 1, \beta = -1, h_{ij} = 0$ and f in (3.6), the assertion of Theorem 3.1 becomes that $A_{s,t}^{f,g}$ is path independent in the sense of (1.7) for some $V \in C^{1,2}(\mathbb{R}_+ \times \mathbb{R}^d; \mathbb{R})$ if and only if

$$\begin{cases} \frac{\partial}{\partial t} V(t, x) = -\frac{1}{2} G\left(\langle \sigma_i, (\nabla^2 V) \sigma_j \rangle\right)(t, x) - \langle \nabla V, b \rangle(t, x), \\ f_{ij}(t, x) = \frac{1}{2} \langle \sigma_i, (\nabla^2 V) \sigma_j \rangle(t, x), \\ g(t, x) = (\sigma^* \nabla V)(t, x), \quad (t, x) \in [0, T] \times \mathbb{R}^d. \end{cases}$$

It is easy to see that this generalizes the main results derived in [16, 23, 24, 26] where $h \equiv 0$ and g is given by $\sigma^{-1}b$, under additional condition ensuring the existence of $\sigma^{-1}b$, i.e., b takes value in $\{\sigma v : v \in \mathbb{R}^d\}$.

However, since G is a linear function in the linear expectation case, Theorem 3.1 does not directly apply to existing results, but extends them to the non-degenerate G -setting.

Moreover, the nonlinear PDE included in (3.1) covers the G -heat equation as a special example.

Remark 3.3. When $h = b = 0$, $\alpha = 1$, and $\beta = -2$, the PDE in (3.1) for V reduces to the following G -heat equation

$$\frac{\partial}{\partial t} V(t, x) + G\left(\left(\langle \sigma_i, (\nabla^2 V) \sigma_j \rangle(t, x)\right)_{1 \leq i, j \leq m}\right) = 0,$$

which is one of main motivations for the study of G -Brownian motion.

4 An Example with $\alpha = 0$

Now we provide an example to demonstrate our main result for $\alpha = 0$. As indicated in Remark 1.1 that when $\alpha \neq 0$ the study can be reduced to $\alpha = 1$.

Example 4.1. Let $d = 1$, $\alpha = 0$, and $\beta = 2$. By Theorem 3.1, $A_{s,t}^{f,g}$ is path independent if and only if

$$(4.1) \quad \begin{cases} f(t, x) = \frac{1}{2} G^{-1} \left(\frac{\partial V}{\partial t} + b \frac{\partial V}{\partial x} \right) (t, x), \\ \left(h \frac{\partial V}{\partial x} \right) (t, x) + \frac{1}{2} \left(\sigma^2 \frac{\partial^2 V}{\partial x^2} \right) (t, x) = 0, \\ g(t, x) = \left(\sigma^* \frac{\partial V}{\partial x} \right) (t, x), \quad (t, x) \in [0, T] \times \mathbb{R}. \end{cases}$$

We may solve V by using \mathcal{L}_t -Harmonic function:

$$(4.2) \quad \mathcal{L}_t V_0(x) = 0, \quad V_0 \in C^{1,2}(\mathbb{R} \rightarrow \mathbb{R}), \quad t \geq 0,$$

where $\mathcal{L}_t = h(t, x) \frac{\partial}{\partial x} + \frac{1}{2} \sigma^2(t, x) \frac{\partial^2}{\partial x^2}$.

For any \mathcal{L}_t -Harmonic function V_0 , $t \geq 0$, let $V(t, x) = \varphi(t) V_0(x)$ for some $\varphi \in C^{1,2}(\mathbb{R}_+ \rightarrow \mathbb{R})$. Then V solves the above PDE in (4.1). Therefore, $A_{s,t}^{f,g}$ is path independent if

$$(4.3) \quad \begin{cases} f(t, x) = \frac{1}{2} G^{-1} \left(\varphi'(t) V_0(x) + b(t, x) \varphi(t) V_0'(x) \right), \\ g(t, x) = \sigma(t, x) \varphi(t) V_0'(x), \quad (t, x) \in [0, T] \times \mathbb{R}. \end{cases}$$

To present specific choices of V_0 , let h and σ do not depend on t . Then (4.2) becomes

$$h(x) V_0'(x) + \frac{1}{2} \sigma^2(x) V_0''(x) = 0.$$

When $\sigma^2(x) \neq 0$, this is equivalent to

$$V_0''(x) = -2 \frac{h(x)}{\sigma^2(x)} V_0'(x).$$

Thus,

$$V_0(x) = V_0(0) + V_0'(0) \int_0^x e^{-2 \int_0^u \frac{h(r)}{\sigma^2(r)} dr} du.$$

In particular, when $\sigma(x) = 1$, $h(x) = x$, we have

$$V_0(x) = V_0(0) + V_0'(0) \int_0^x e^{-u^2} du,$$

which is related to the Gaussian distribution.

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