

Chaotic Kabanov formula for the Azéma martingales

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Abstract

We derive the chaotic expansion of the product of n -th and first order multiple stochastic integrals with respect to certain normal martingales. This is done by application of the classical and quantum product formulas for multiple stochastic integrals. Our approach extends existing results on chaotic calculus for normal martingales and exhibits properties relative to multiple stochastic integrals, polynomials and Wick products, that characterize the Wiener and Poisson processes.

Key words: Azéma martingales, multiple stochastic integrals, product formulas.

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1 Introduction

The Wiener-Itô and Poisson-Itô chaotic decompositions give an isometric isomorphism between the Fock space $\Gamma(L^2(\mathbb{R}_+))$ and the space of square-integrable functionals of the process. This isomorphism is constructed by association of a symmetric function $f_n \in L^2(\mathbb{R}_+)^{\circ n}$ to its multiple stochastic integral. The isometry property comes from the fact that the angle brackets $\langle B \rangle_t$ and $\langle \tilde{N} \rangle_t$ of both the Wiener and compensated Poisson processes $(B_t)_{t \in \mathbb{R}_+}$, $(\tilde{N}_t)_{t \in \mathbb{R}_+}$ are equal to t . Therefore such isometric isomorphisms may exist for more general processes, provided their angle bracket equals t .

A martingale is said to be normal (Dellacherie and Meyer, 1992) if its angle bracket is equal to t and if the range of the multiple stochastic integrals isometry defined on $\Gamma(L^2(\mathbb{R}_+))$ is equal to the space of square-integrable functionals of this martingale. This last property is called the chaotic representation property, or CRP.

The quadratic variation $([M, M]_t)_{t \in \mathbb{R}_+}$ of a normal martingale $(M_t)_{t \in \mathbb{R}_+}$ in $L^4(\Omega)$ with the CRP (or less restrictively, with the predictable representation property) satisfies the equation

$$[M, M]_t = t + \int_0^t \phi_s dM_s, \quad t \in \mathbb{R}_+, \quad (1)$$

called a structure equation (Emery, 1990), where $(\phi_s)_{s \in \mathbb{R}_+}$ is a predictable process. On the other hand, a martingale satisfying (1) has angle bracket equal to t , but does not necessarily possess the chaotic representation property.

In this paper we work with solutions of such equations, which include the Wiener process, $\phi = 0$, the compensated Poisson processes, ϕ constant, nonzero, and the Azéma martingales, $\phi_t = \beta M_t$, $\beta \in [-2, 0]$.

Note that in the probabilistic framework of the structure equation one has to use a predictable version of ϕ , i.e. $\phi_t = \beta M_{t-}$ in the case of the Azéma martingale. However, since we are working in L^2 with stochastic integrals with respect to a normal martingale, it follows from a remark of Dellacherie and Meyer (1992, p. 199) that we need not distinguish between the adapted and predictable versions of ϕ . Thus, in the sequel we will use $\phi_t = \beta M_t$ in the case of the Azéma martingale.

In the stochastic analysis on both the Wiener and Poisson spaces, multiplication formulas for multiple stochastic integrals proved to be useful tools as they give the chaotic expansion of the product of two multiple stochastic integrals. Recently (Russo and Vallois, 1998), multiplication formulas have been proposed for multiple stochastic integrals with respect to normal martingales. However, these formulas do not give the explicit chaotic expansion of the considered products.

Our first goal in this paper is to relate the product formulas of Kabanov (1975), Russo and Vallois (1998) and Surgailis (1984) to their counterparts in quantum probability, cf. Section 3. With these tools we compute the chaotic expansion of the product of a multiple stochastic integral with a single stochastic integral, and refer to this formula as the chaotic Kabanov formula. This formula is proved for a class of martingales that includes the Azéma martingales. We have not obtained a

general formula for the chaotic expansion of the product of two multiple stochastic integrals of orders n and m . This seems to require complicated computations.

As an application of this Kabanov formula we show that a number of properties in stochastic analysis such as the possibility of expressing multiple stochastic integrals with polynomials and some properties of the Wick product are specific to the Wiener and Poisson cases, cf. Section 5. In the last section we consider formulas giving the derivation of a product.

2 Notation and preliminaries

We denote by

$$\Gamma(L^2(\mathbb{R}_+)) = \bigoplus_{n \in \mathbb{N}} L^2(\mathbb{R}_+)^{\otimes n}$$

the symmetric Fock space over $L^2(\mathbb{R}_+)$, where $L^2(\mathbb{R}_+)^{\otimes n}$ is endowed with the norm $\|\cdot\|_{L^2(\mathbb{R}_+)^{\otimes n}} = n! \|\cdot\|_{L^2(\mathbb{R}_+)^{\otimes n}}$, and $\langle \cdot, \cdot \rangle$ denotes the scalar product in $L^2(\mathbb{R}_+)$.

In this paper we work with a normal martingale $(M_t)_{t \in \mathbb{R}_+}$, with $M_0 = 0$, that satisfies the structure equation (1).

The multiple stochastic integral of a symmetric function $f_n \in L^2(\mathbb{R}_+)^{\otimes n}$ with respect to $(M_t)_{t \in \mathbb{R}_+}$ is denoted by

$$I_n(f_n) = n! \int_0^\infty \int_0^{t_1^-} \cdots \int_0^{t_{n-1}^-} f_n(t_1, \dots, t_n) dM_{t_1} \cdots dM_{t_n},$$

and as a convention we let

$$I_n(f_n) = I_n(\tilde{f}_n) = I_n(f_n(t_1, \dots, t_n)), \quad (2)$$

if $f_n \in L^2(\mathbb{R}_+)^{\otimes n}$ is not symmetric, where \tilde{f}_n denotes the symmetrization of f_n in its n variables.

We denote by \mathcal{F}_t the σ -algebra generated by $\{M_s : 0 \leq s \leq t\}$, $t \in \mathbb{R}_+$.

We let $\nabla^- : \Gamma(L^2(\mathbb{R}_+)) \rightarrow \Gamma(L^2(\mathbb{R}_+)) \otimes L^2(\mathbb{R}_+)$ and $\nabla^+ : \Gamma(L^2(\mathbb{R}_+)) \otimes L^2(\mathbb{R}_+) \rightarrow \Gamma(L^2(\mathbb{R}_+))$ denote the annihilation and creation operators on $\Gamma(L^2(\mathbb{R}_+))$, defined as

$$\nabla_t^- I_n(f_n) = n I_{n-1}(f_n(\cdot, t)), \quad t \in \mathbb{R}_+, \quad \nabla^+(I_n(f_{n+1})) = I_{n+1}(f_{n+1}),$$

where f_{n+1} is symmetric in its n first variables, $n \in \mathbb{N}$.

Let $\mathcal{A} = \{]a, b] : a, b \in \mathbb{R}_+, a < b\}$. For $]a, b],]c, d] \in \mathcal{A}$ we use the notation $]a, b] \leq]c, d]$ if $b \leq c$. Let \mathcal{S}_n , $n \in \mathbb{N}$, denote the vector space

$$\mathcal{S}_n = \left\{ \sum_{i=1}^{i=m} \alpha_i I_n(1_{A_1^i} \circ \cdots \circ 1_{A_n^i}) : A_k^i \cap A_l^i = \emptyset, 1 \leq k < l \leq n, A_1^i, \dots, A_n^i \in \mathcal{A}, \alpha_i \in \mathbb{R}, i = 1, \dots, m, m \geq 1 \right\},$$

and let \mathcal{S} denote the vector space generated by $\cup_{n \in \mathbb{N}} \mathcal{S}_n$, which is dense in $L^2(\Omega)$.

Let \mathcal{V}_n denote the vector space

$$\mathcal{V}_n = \left\{ \sum_{i=1}^{i=m} \alpha_i 1_{A_0^i} I_n(1_{A_1^i} \circ \cdots \circ 1_{A_n^i}) : A_k^i \cap A_l^i = \emptyset, 1 \leq k < l \leq n, A_0^k \cap A_0^l = \emptyset, 1 \leq k < l \leq m, A_0^i, \dots, A_n^i \in \mathcal{A}, \alpha_i \in \mathbb{R}, i = 1, \dots, m, m \geq 1 \right\}.$$

We denote by \mathcal{V} the vector space generated by $\cup_{n \geq 0} \mathcal{V}_n$, which is dense in $L^2(d\lambda)$ and in $L^1(d\lambda)$, where $d\lambda = dt \times dP$.

Let \mathcal{U}_n denote the vector space

$$\mathcal{U}_n = \left\{ \sum_{i=1}^{i=m} \alpha_i 1_{A_0^i} I_n(1_{A_1^i} \circ \cdots \circ 1_{A_n^i}) : A_k^i \cap A_l^i = \emptyset, 0 \leq k < l \leq n, A_0^k \cap A_0^l = \emptyset, 1 \leq k < l \leq m, A_0^i, \dots, A_n^i \in \mathcal{A}, \alpha_i \in \mathbb{R}, i = 1, \dots, m, m \geq 1 \right\}.$$

Note that \mathcal{U}_n is made of processes in \mathcal{V}_n that do not depend on the ‘‘present’’.

We denote by \mathcal{U} the vector space generated by $\cup_{n \geq 0} \mathcal{U}_n$, and by $\bar{\mathcal{U}}$ the completion of \mathcal{U} in $L^2(d\lambda)$. The space $\bar{\mathcal{U}}$ is strictly smaller than $L^2(d\lambda)$ but it contains the adapted square-integrable processes and the range of ∇^- , since $F \in \mathcal{S}$ implies $\nabla^- F \in \mathcal{U}$.

We will use the following lemma of Ma, Protter and San Martin (1998).

Lemma 1 *Let $I_n(f_n), I_1(g_1) \in \mathcal{S}$. We have*

$$I_n(f_n)I_1(g_1) = I_{n+1}(f_n \circ g_1) + n \int_0^\infty I_{n-1}(f_n(\cdot, t))g_1(t)d[M, M]_t. \quad (3)$$

As a consequence of this lemma, we have

$$I_n(1_{A_1} \circ \cdots \circ 1_{A_n}) = I_1(1_{A_1}) \cdots I_1(1_{A_n})$$

whenever $A_i \cap A_j = \emptyset$, $1 \leq i < j \leq n$.

We assume that there exists a set \mathcal{P} of functionals dense in $L^2(\Omega)$ and included both in $L^\infty(\Omega)$ and in $Dom(\nabla^-)$, and such that \mathcal{P} is stable by ∇_t^- , $t \in \mathbb{R}_+$. This assumption is satisfied in the cases of interest to us, that is in the Wiener and Poisson cases, and also in the case of the Azéma martingales since the later is bounded (in this case it suffices to take $\mathcal{P} = \mathcal{S}$).

We recall the following identity, satisfied in general on Fock space:

$$E[\nabla^+(u)^2] = \|u\|_{L^2(\Omega) \otimes L^2(\mathbb{R}_+)}^2 + E \left[\int_0^\infty \int_0^\infty \nabla_s^- u_t \nabla_t^- u_s ds dt \right]. \quad (4)$$

Definition 1 Let $Dom_{1,2}(\nabla^+)$ denote the set of processes $u \in L^1(d\lambda)$ such that there exists a sequence $(u_n)_{n \in \mathbb{N}} \subset \mathcal{V}$ converging in $L^1(d\lambda)$ to u and such that $(\nabla^+(u_n))_{n \in \mathbb{N}}$ converges in $L^1(\Omega)$ to an element of $L^2(\Omega)$ denoted as $\nabla^+(u)$.

Remark 1 The domain $Dom_{1,2}(\nabla^+)$ is well-defined.

Proof. Assume that $(u_n)_{n \in \mathbb{N}} \subset \mathcal{V}$ converges in $L^1(d\lambda)$ to 0 and that $(\nabla^+(u_n))_{n \in \mathbb{N}}$ converges in $L^1(\Omega)$ to $G \in L^2(\Omega)$. From the duality relation

$$E[\nabla^+(u_n)F] = E \left[\int_0^\infty u_n(t) \nabla_t^- F dt \right], \quad F \in \mathcal{P},$$

we obtain $E[GF] = 0$, $F \in \mathcal{P}$, hence $G = 0$ since $G \in L^2(\Omega)$. □

This also means that if $u \in Dom_{1,2}(\nabla^+)$ and there exists a sequence $(u_n)_{n \in \mathbb{N}}$ converging in $L^1(\Omega)$ to u and such that $(\nabla^+(u_n))_{n \in \mathbb{N}}$ converges in $L^1(\Omega)$ to $G \in L^2(\Omega)$, then $G = \nabla^+(u)$. Naturally, $Dom_{1,2}(\nabla^+)$ contains the usual L^2 domain of ∇^+ which is denoted as $Dom_2(\nabla^+)$.

Finally, let

$$\mathcal{T}_n = \left\{ \sum_{i=1}^{i=m} F_i 1_{]a_i, b_i]} I_n(1_{A_1^i} \circ \dots \circ 1_{A_n^i}) : [0, b_i] \cap A_j^i = \emptyset, j = 1, \dots, n, \right. \\ \left. A_k^i \cap A_l^i = \emptyset, 1 \leq k < l \leq n, A_1^i, \dots, A_n^i \in \mathcal{A}, F_i \in L^2(\mathcal{F}_{a_i}), i = 1, \dots, m, m \geq 1 \right\},$$

and \mathcal{T} the vector space generated by $\cup_n \mathcal{T}_n$.

The space $\mathcal{T} \cap L^2(d\lambda)$ is not dense in $L^2(d\lambda)$, but it contains \mathcal{U} , and it is dense in $\bar{\mathcal{U}}$.

Remark 2 Let $u, v \in \mathcal{U}$, where v is an adapted process. Then $uv \in \mathcal{T}$.

Proof. We write $u = 1_A I_n(1_{A_1} \circ \cdots \circ 1_{A_n})$ with $A_1, \dots, A_k \leq B$ and

$$v = 1_B I_m(1_{A_{n+1}} \circ \cdots \circ 1_{A_{n+m}})$$

with $A_{n+1}, \dots, A_{n+m} \leq B$ since v is adapted. Then

$$uv = 1_{A \cap B} I_{k+m}(1_{A_1} \circ \cdots \circ 1_{A_k} \circ 1_{A_{n+1}} \circ \cdots \circ 1_{A_{n+m}}) I_{n-k}(1_{A_{k+1}} \circ \cdots \circ 1_{A_n}),$$

hence $uv \in \mathcal{T}$. □

Proposition 1 For $u \in \mathcal{T}$ such that $\int_0^\infty u_t dM_t \in L^2(\Omega)$, we have $u \in \text{Dom}_{1,2}(\nabla^+)$ and

$$\nabla^+(u) = \int_0^\infty u_t dM_t. \quad (5)$$

Proof. We start by choosing u of the form $u = F 1_{]a,b]} I_n(1_{A_1} \circ \cdots \circ 1_{A_n})$, with $F = I_m(1_{B_1} \circ \cdots \circ 1_{B_m})$, $A_i \cap [0, b] = \emptyset$, $i = 1, \dots, n$, and $B_1, \dots, B_m \subset [0, a]$. Then from Lemma 1:

$$\begin{aligned} \nabla^+(u) &= I_{n+m+1}(1_{B_1} \circ \cdots \circ 1_{B_m} \circ 1_{]a,b]} \circ \cdots \circ 1_{A_n}) \\ &= I_1(1_{]a,b]}) I_m(1_{B_1} \circ \cdots \circ 1_{B_m}) I_n(1_{A_1} \circ \cdots \circ 1_{A_n}) \\ &= F I_1(1_{]a,b]}) I_n(1_{A_1} \circ \cdots \circ 1_{A_n}). \end{aligned}$$

Hence, by linearity, for any $u \in \mathcal{T}$ and $F \in \mathcal{S}$,

$$\nabla^+(u) = F I_1(1_{]a,b]}) I_n(1_{A_1} \circ \cdots \circ 1_{A_n}) = \int_0^\infty u_t dM_t.$$

If $u \in \mathcal{T}$, then choosing a sequence $(F_n)_{n \in \mathbb{N}} \subset \mathcal{S}$ that converges in $L^2(\Omega)$ to F , and letting $u_m = F_m 1_{]a,b]} I_n(1_{A_1} \circ \cdots \circ 1_{A_n})$, we have the convergence of $(u_m)_{m \in \mathbb{N}}$ in $L^1(d\lambda)$ to u and of $(\nabla^+(u_m))_{m \in \mathbb{N}}$ in $L^1(\Omega)$ to $\int_0^\infty u_t dM_t$ which belongs to $L^2(\Omega)$ by hypothesis. □

3 Classical and quantum product formulas

The aim of this section is to link two versions of the product formula for a multiple stochastic integral and a single stochastic integral with respect to normal martingales.

The first one is called the Kabanov formula in reference to Kabanov (1975), which treated the Poisson case, i.e. $\phi = 1$, and can be stated as follows. For $f_n \in L^2(\mathbb{R}_+)^{\circ n}$ and $g \in L^2([0, T])$, in particular g with compact support, Lemma 1 was extended by Ma, Protter and San Martin(1998) and Russo and Vallois (1998) as

$$I_n(f_n)I_1(g) = I_{n+1}(f_n \circ g) + n \int_0^\infty I_{n-1}^*(f_n(\cdot, t))g(t)d[M, M]_t. \quad (6)$$

The second term in this formula is an integration over a diagonal, due to the Itô formula, and the notation “ $I_{n-1}^*(f_n(\cdot, t))$ ” will be made precise in Def. 2.

The second product formula uses only chaos expansions and Fock space, and can be found in the work of quantum probabilists. It is often stated in the formalism of *quantum stochastic integrals*:

$$\int_0^\infty g_t dM_t = \int_0^\infty g_t da_t^+ + \int_0^\infty g_t da_t^- + \int_0^\infty \phi_t g_t da_t^\circ, \quad (7)$$

cf. (Parthasarathy, 1990) and (Biane, 1993, Proposition 18).

The definition of quantum stochastic integrals as operators poses several functional analytic problems. See the paper of Attal and Lindsay (1997) for recent extensions of their definition. In general the above relation does not hold in $L^2(\Omega)$ but only in a weak sense (Biane, 1993), where this formula is proved for F an exponential vector, and for bounded predictable g and ϕ , with compact support.

Relation (7) can be reformulated formally using the operators ∇^- , ∇^+ , as follows:

$$F\nabla^+(g) = \nabla^+(F \otimes g) + \int_0^\infty g(t)\nabla_t^- F dt + \nabla^+(g\phi\nabla^- F), \quad (8)$$

cf. e.g. (Attal, 1997), (Biane, 1993), and the references therein, the only difference between (7) and (8) being a change of notation.

In the Poisson case, this formula also appears in the papers of Nualart and Vives (1990) and Dermoune, Krée and Wu (1988). It will play an important role in the computation of the chaotic expansion of the Kabanov formula, since unlike (3) and (6) it uniquely involves calculations on chaos.

The product formula (8) can be rewritten in the language of quantum stochastic differentials as

$$\int_0^\infty g(t) dM_t = \int_0^\infty g(s) dB_s + \int_0^\infty \phi(s) g(s) da_t^\circ,$$

or

$$dM_t = dB_t + \phi(t) da_t^\circ. \tag{9}$$

In the Wiener interpretation of the Fock space, the operator differential dB_t identifies to the multiplication by the classical Brownian differential, and this formula states that the operator differential $dM_t = dB_t + \phi(t) da_t^\circ$ is identified to the multiplication operator by the classical differential dM_t when the Fock space is identified to the L^2 space of $(M_t)_{t \in \mathbf{R}_+}$. However, the equation (9) has no classical interpretation because the operator processes $(B_t)_{t \in \mathbf{R}_+}$ and $(a_t^\circ)_{t \in \mathbf{R}_+}$ can not be interpreted simultaneously as multiplication operators in the same probabilistic interpretation of the Fock space (the reason for this is that they do not commute). Consequently, (9) does not have a classical meaning; it only defines an operator process *on Fock space*.

If $\phi(t)$ is a function of M_t , e.g. $\phi(t) = f(M_t)$, then (9) becomes a *quantum stochastic differential equation*, in the space of operators on Fock space:

$$dM_t = dB_t + f(M_t) da_t^\circ, \tag{10}$$

that does not have a classical interpretation, whatever the interpretation chosen for the Fock space.

In order to find the multiplication formula for multiple stochastic integrals with respect to $(M_t)_{t \in \mathbf{R}_+}$, one has to compute in particular the chaos expansion of M_t multiplied by a multiple stochastic integral. This means that the explicit expression of M_t *as an operator on Fock space* has to be obtained. A way to obtain this expression is to solve (10) *in the space of operators on Fock space*, that is to determine the process of operators associated to $(M_t)_{t \in \mathbf{R}_+}$. The simplest case is the

linear case, i.e. $f(x) = \beta x$, $\beta \in [-2, 0]$, and it corresponds to the family of Azéma martingales. In this case, (10) reads

$$dM_t = dB_t + \beta M_t da_t^\circ, \quad (11)$$

and it can be formally solved as a linear equation. It can be easily shown that its solution is

$$M_t = \sum_{n \geq 0} \beta^n \int_0^t \int_0^{t_{n-1}} \cdots \int_0^{t_2} B_{t_1} da_{t_2}^\circ \cdots da_{t_{n-1}}^\circ,$$

or

$$M_t = \sum_{n \geq 0} \beta^n \int_0^t \int_0^{t_n} \cdots \int_0^{t_1} dB_{t_0} da_{t_1}^\circ \cdots da_{t_n}^\circ.$$

This expression can now be rewritten explicitly in terms of operators on Fock space:

$$\begin{aligned} M_t F &= \sum_{n \geq 0} \beta^n \nabla^+ \left(1_{[0,t]} \left(\nabla^+ \left(\cdots \nabla^+ \left(\int_0^\cdot \nabla_s^- ds \nabla^- \cdots \nabla^- F \right) \right) \right) \right) \quad (12) \\ &+ \sum_{n \geq 0} \beta^n \nabla^+ (1_{[0,t]} (\nabla^+ (\cdots \nabla^+ (\nabla^+ (1_{[0,\cdot]} \nabla^- \cdots \nabla^- F))))). \end{aligned}$$

When $F = I_n(f_n)$, this gives the explicit chaos expansion of $M_t I_n(f_n)$.

The above calculation is formal and some further computations are required in order to obtain the explicit chaos expansion of $M_t I_n(f_n)$. In this paper we justify each of the above steps by explicit calculations and carry out the final computations suggested in (12). The aim of this remark was to show that the quantum stochastic point of view gives a different understanding of the problem and provides a quick solution by reducing the problem to the determination of a quantum diffusion. We stress that although the solution process has a classical version, the diffusion equation (10) is meaningful only in the space of operators on Fock space. Also, those remarks show that the reason why explicit calculation can be carried out is in fact that (11) is a *linear* equation.

We also note that a slightly more general case can be considered, i.e. $\phi(t) = \lambda_t + \int_0^t \beta_t(s) dM_s$, which is still linear. In this case the solution reads

$$\begin{aligned} M_t &= \sum_{n \geq 0} \int_0^t \int_0^{t_n} \beta_{t_n}(t_{n-1}) \int_0^{t_{n-1}} \cdots \int_0^{t_1} dB_{t_0} da_{t_1}^\circ \cdots da_{t_n}^\circ \\ &+ \sum_{n \geq 0} \int_0^t \int_0^{t_n} \beta_{t_n}(t_{n-1}) \int_0^{t_{n-1}} \cdots \int_0^{t_2} \beta_{t_2}(t_1) \lambda_{t_1} da_{t_0}^\circ \cdots da_{t_n}^\circ. \end{aligned}$$

This situation is considered in our paper.

Our aim in this section is to prove (7) in $L^2(\Omega)$ from (6). The proof of Biane (1993), cf. also (Attal, 1997), relies on the construction of quantum stochastic calculus, whereas our proof uses more classical probabilistic arguments. We also prove that (8) holds in L^2 under assumptions that are satisfied in our setting.

Definition 2 *Following Ma, Protter and San Martin (1998) and Russo and Vallois (1998), let μ denote the measure on $\mathbb{R}_+ \times \Omega$ defined as $\mu([0, t] \times A) = E[1_A[M, M]_t]$, $A \in \mathcal{F}$, $t \in \mathbb{R}_+$ and let $\nu = (\lambda + \mu)/2$. For $u \in \bar{\mathcal{U}}$ we denote by u^* the limit in $L^2(d\nu)$ of any sequence $(u_n)_{n \in \mathbb{N}} \subset \mathcal{U}$ that converges in $L^2(d\lambda)$ to u .*

Note that $u = u^*$ λ -a.e., but not μ -a.e., except in the Wiener case ($\phi = 0$), and in general for $u \in \mathcal{U}$. From Lemma 5.2. of Ma, Protter and San Martin (1998) we have

$$\| u \|_{L^2(d\lambda)} = \| u^* \|_{L^2(d\mu)}, \quad u \in \bar{\mathcal{U}}.$$

In what follows we fix $T > 0$. For $u \in \bar{\mathcal{U}}$, $\int_0^T u_t^* d[M, M]_t$ is defined a.e. as an integral with respect to an increasing process and it belongs to $L^1(\Omega)$, since

$$E \left[\left| \int_0^T u_t^* d[M, M]_t \right| \right] \leq T^{\frac{1}{2}} \| u^* \|_{L^2(d\mu)} \leq T^{\frac{1}{2}} \| u \|_{L^2(d\lambda)}. \quad (13)$$

The following Lemma allows us to prove the quantum product formula (7) from the probabilistic product formula (6).

Lemma 2 *Let $u \in \bar{\mathcal{U}}$ and assume that $1_{[0, T]} \phi \in L^2(d\lambda)$ for some $T \in \mathbb{R}_+$. Then $1_{[0, T]} u \phi \in \text{Dom}_{1,2}(\nabla^+)$ if and only if*

$$\int_0^T u_t^* d[M, M]_t \in L^2(\Omega),$$

and in this case,

$$\int_0^T u_t dt + \nabla^+(1_{[0, T]} u \phi) = \int_0^T u_t^* d[M, M]_t. \quad (14)$$

Proof. We start by assuming that $u \in \mathcal{U}$ is of the form $u = \sum_{i=1}^n F_i 1_{[a_i, b_i]}$ where $b_i \leq T$. Since ϕ is adapted we can choose a sequence $(v_n)_{n \in \mathbb{N}} \subset \mathcal{U}$ of adapted processes converging to $\phi 1_{[0, T]}$ in $L^2(d\lambda)$. Then, from Remark 2, $u v_n \in \mathcal{T}$ and from Proposition 1,

$$(\nabla^+(uv_n))_{n \in \mathbb{N}} = \left(\int_0^\infty uv_n(t) dM_t \right)_{n \in \mathbb{N}}.$$

Hence $(\nabla^+(uv_n))_{n \in \mathbb{N}}$ converges in $L^1(\Omega)$ to

$$\sum_{i=1}^{i=n} F_i \int_{a_i}^{b_i} \phi_s dM_s = \int_0^T u_t \phi_t dM_t \in L^2(\Omega),$$

$u\phi \in Dom_{1,2}(\nabla^+)$ and $\nabla^+(u\phi) = \int_0^T u_t \phi_t dM_t$, i.e.

$$\int_0^T u_t dt + \nabla^+(u\phi) = \int_0^T u_t d[M, M]_t.$$

If $u \in \bar{\mathcal{U}}$ we approximate it in $L^2(d\lambda)$ by a sequence $(u_n)_{n \in \mathbb{N}}$ in \mathcal{U} . Then, from (13) we have that, $\left(\int_0^T u_n(t) dt + \nabla^+(u_n\phi) \right)_{n \in \mathbb{N}} = \left(\int_0^T u_n(t) d[M, M]_t \right)_{n \in \mathbb{N}}$ converges in $L^1(\Omega)$ to $\int_0^T u^* d[M, M]_t$. □

We now prove the quantum product formula (7), under assumptions different from that of Biane (1993).

Proposition 2 *Let F be in a finite sum of chaos and let $h \in L^2([0, T])$ be bounded. Then $h\phi\nabla^- F \in Dom_{1,2}(\nabla^+)$ if and only if $I_1(h)F \in L^2(\Omega)$ and in this case,*

$$I_1(h)F = \nabla^+(h \otimes F) + (h, \nabla^- F)_{L^2(\mathbb{R}_+)} + \nabla^+(h\phi\nabla^- F). \quad (15)$$

Proof. We let $F = I_n(f_n) \in \mathcal{S}$, $h \in L^2([0, T])$, and apply (6) and the above lemma to $h\nabla^- I_n(f_n) \in \bar{\mathcal{U}}$, since $\bar{\mathcal{U}}$ contains the range of ∇^- . In the general case we choose a sequence $(F_n)_{n \in \mathbb{N}} \subset \mathcal{S}$ that converges to F in $L^2(\Omega)$. The right-hand side of (15) converges to $\nabla^+(h \otimes F) + \int_0^T h_t \nabla_t^- F dt + \nabla^+(h\phi\nabla^- F) = I_1(h)F$ in $L^1(\Omega)$ and $(h\phi\nabla^- F_n)_{n \in \mathbb{N}}$ converges to $h\phi\nabla^- F$ in $L^1(\Omega)$ as n goes to infinity. Hence $h\phi\nabla^- F \in Dom_{1,2}(\nabla^+)$ if and only if $I_1(h)F \in L^2(\Omega)$. □

The result extends to bounded simple adapted processes of the form $h = G1_{[t_1, t_2[}$ for G , \mathcal{F}_{t_1} -measurable, since in this case $\nabla_s^- G = 0$, $s > t_1$, cf. Lemma 4.1. of Ma, Protter and San Martin (1998), which means that $\nabla^+(1_{[t_1, t_2[} G) = G\nabla^+(1_{[t_1, t_2[})$.

We close this section with a remark on the link between independence of stochastic integrals and their deterministic kernels. The above formula easily gives

information on the chaotic expansion of a product of stochastic integrals, for general ϕ . If $f, g \in L^2([0, T])$ are such that $fg\phi \in Dom_{1,2}(\nabla^+)$, then

$$I_1(f)I_1(g) = I_2(f \circ g) + (f, g)_{L^2(\mathbb{R}_+)} + \nabla^+(fg\phi). \quad (16)$$

Letting T go to infinity, the formula holds for $f, g \in L^2(\mathbb{R}_+)$, provided $fg\phi \in Dom_{1,2}(\nabla^+)$. From this formula it is clearly seen that the chaotic expansion of $I_1(f)I_1(g)$ may be an infinite sum of multiple stochastic integrals, depending on the chaotic expansion of ϕ . This formula can be applied in order to obtain a necessary condition for the independence of stochastic integrals. The result of Urbanik (1967) says that if a stochastic process $(X_t)_{t \in \mathbb{R}_+}$ has stationary and independent increments then independence of $\int_0^\infty f_t dX_t$ and $\int_0^\infty g_t dX_t$ implies $fg = 0$, except if X is Gaussian, in which case the condition becomes $(f, g)_{L^2(\mathbb{R}_+)} = 0$. In the Gaussian case this property has been further extended to multiple stochastic integrals by Üstünel and Zakai (1990). See also Privault (1996) for the case of deterministic ϕ . We can now give an extension of this property to more general ϕ .

Proposition 3 *Assume that the chaotic expansion of ϕ_t does not contain terms of order one. Then, independence of $I_1(f)$ and $I_1(g)$ implies $(f, g)_{L^2(\mathbb{R}_+)} = 0$ and $f_t g_t = 0$, $E[\phi_t^2] dt$ -a.e.*

Proof. We follow an argument of Üstünel and Zakai (1990) that deals with the Wiener case.

$$\begin{aligned} \|f \circ g\|_{L^2(\mathbb{R}_+)^{\circ 2}}^2 &\geq (f, f)_{L^2(\mathbb{R}_+)}(g, g)_{L^2(\mathbb{R}_+)} = E[I_1(f)^2]E[I_1(g)^2] \\ &= E[(I_1(f)I_1(g))^2] \\ &= \|f \circ g\|_{L^2(\mathbb{R}_+)^{\circ 2}}^2 + (f, g)_{L^2(\mathbb{R}_+)}^2 + \|\nabla^+(fg\phi)\|_{L^2(\Omega)}^2, \end{aligned}$$

where we used the assumption on the chaotic expansion of ϕ which implies the orthogonality of $\nabla^+(fg\phi)$ and $I_2(f \circ g)$ in (16). Hence $(f, g)_{L^2(\mathbb{R}_+)} = 0$, and $\nabla^+(fg\phi) = 0$ a.e. Hence $fg\phi = 0$, $d\lambda$ -a.e. □

Note that the hypothesis of this proposition do not include the Azéma martingales.

4 Chaotic Kabanov formula

The purpose of this section is to obtain in Theorem 1 a chaotic formula for the product $I_1(g)I_n(f_n)$ using Proposition 2 applied to $F = I_n(f_n)$. Under certain assumptions on f_n and g , we have

$$I_1(g)I_n(f_n) = I_{n+1}(f_n \circ g) + nI_{n-1}(\langle f_n(*, \cdot), g(\cdot) \rangle) + n\nabla^+(g(\cdot)\phi(\cdot)I_{n-1}(f_n(*, \cdot))). \quad (17)$$

If ϕ is deterministic, in particular in the Wiener and Poisson cases, this formula easily yields the Kabanov formula (Privault, 1996). On the other hand, if ϕ is random the chaos expansion of the term $\phi I_{n-1}(f_n(*, \cdot))$ is unknown unless ϕ_t belongs to the first chaos, i.e. $\phi(t) = \alpha(t) + I_1(\beta_t)$, $t \geq 0$, since in this case an induction argument can be used in (17) to compute the term $\phi I_{n-1}(f_n(t_1, \dots, t_{n-1}, \cdot))$ and to determine the chaotic expansion of the product $I_1(\beta_t)I_n(f_n)$. Thus in this case the chaotic expansion of $I_1(g)I_n(f_n)$ can be obtained as a consequence of Proposition 2.

In this section, we assume that ϕ is of the form $\phi(t) = \alpha(t) + I_1(\beta_t)$, $t \geq 0$, where α is a locally bounded function from \mathbb{R}_+ to \mathbb{R} and $\beta_t(\cdot)$ is a bounded function from \mathbb{R}_+ to \mathbb{R} with support in $[0, t]$, for any fixed $t \in \mathbb{R}_+$. This situation is more general than the situation of Russo and Vallois (1996, Section 4); moreover chaotic expansions are completely determined here.

We have the following lemma:

Lemma 3 *Let $\theta \in \mathbb{R}_+$. Let $f_p \in L^2(\mathbb{R}_+)^{\circ p}$ and*

$$\zeta(t_1, \dots, t_{k-1}; t_k) = \beta_{t_2}(t_1) \cdots \beta_{t_k}(t_{k-1}), \quad t_1, \dots, t_k \in \mathbb{R}_+, \quad k > 1,$$

and $\zeta(t) = 1$, $t \in \mathbb{R}_+$, for $k = 1$. The chaos expansion of the product $I_1(\beta_\theta)I_p(f_p)$ is given by

$$I_1(\beta_\theta)I_p(f_p) = I_{p+1}(g_{p+1}^+(*; \theta)) + I_p(g_p^\circ(*; \theta)) + I_{p-1}(g_{p-1}^-(*; \theta)),$$

where

$$g_{p+1}^+(t_1, \dots, t_{p+1}; \theta) = \sum_{j=1}^{p+1} \frac{p!}{(j-1)!} \zeta(t_j, \dots, t_{p+1}; \theta) f_p(t_1, \dots, \hat{t}_j, \dots, t_{p+1}),$$

where \hat{t}_j means that t_j is omitted in the arguments of f_p ,

$$g_p^\circ(t_1, \dots, t_p; \theta) = f_p(t_1, \dots, t_p) \sum_{j=1}^p \frac{p!}{(j-1)!} \zeta(t_j, \dots, t_p; \theta) \alpha(t_j),$$

and

$$g_{p-1}^-(t_1, \dots, t_{p-1}; \theta) = \sum_{j=1}^p \frac{p!}{(j-1)!} \zeta(t_j, \dots, t_{p-1}; \theta) \langle f_p(t_1, \dots, t_{p-1}, \cdot), \beta_{t_j}(\cdot) \rangle,$$

with $t_p = \theta$.

Proof. We will prove the lemma by induction on p . Observe that the kernels g^+ , g° and g^- are not symmetric functions and that this result uses the convention (2).

1. For $p = 1$, using Proposition 2, we have

$$\begin{aligned} I_1(\beta_\theta)I_1(f) &= I_2(\beta_\theta \circ f) + \langle \beta_\theta, f \rangle + \nabla^+(\beta_\theta \phi f) \\ &= I_2(\beta_\theta \circ f) + \langle \beta_\theta, f \rangle + I_1(\beta_\theta \alpha f) + I_2(\beta_{\theta(t_1)} f(t_1) \circ \beta_{t_1}(t_2)) \\ &= I_2(\beta_\theta \circ f + \beta_\theta f \circ \beta_t) + I_1(\beta_\theta \alpha f) + \langle \beta_\theta, f \rangle, \end{aligned}$$

and therefore, in particular, $I_1(\beta_\theta)I_1(f)$ is in the domain of ∇^+ .

2. For the general case, applying Proposition 2 and the fact that f_{p+1} is symmetric, we have

$$\begin{aligned} I_1(\beta_\theta)I_{p+1}(f_{p+1}) &= I_{p+2}(\beta_\theta \circ f_{p+1}) + (p+1)I_p(\langle \beta_\theta(\cdot), f_{p+1}(*, \cdot) \rangle) \\ &\quad + (p+1)\nabla^+(\beta_\theta(\cdot)\alpha(\cdot)I_p(f_{p+1}(*, \cdot))) \\ &\quad + (p+1)\nabla^+(\beta_\theta(\cdot)I_1(\beta_\cdot)I_p(f_{p+1}(*, \cdot))). \end{aligned}$$

Now, applying the induction hypothesis to $I_1(\beta_{t_{p+2}})I_p(f_{p+1}(*, t_{p+2}))$ into the last term, we have,

$$\begin{aligned}
I_1(\beta_\theta)I_{p+1}(f_{p+1}) &= I_{p+2}(\beta_\theta \circ f_{p+1}) \\
&+ (p+1)I_{p+2}(\beta_\theta(t_{p+2}) \circ \eta_{p+1}^+(t_1, \dots, t_{p+1}; t_{p+2})) \\
&+ (p+1)I_{p+1}(\beta_\theta(t_{p+1}) \alpha(t_{p+1}) f_{p+1}(t_1, \dots, t_{p+1})) \\
&+ (p+1)I_{p+1}(\beta_\theta(t_{p+2}) \circ \eta_p^0(t_1, \dots, t_p; t_{p+2})) \\
&+ (p+1)I_p(\langle \beta_\theta(\cdot), f_{p+1}(\cdot, \cdot) \rangle) \\
&+ (p+1)I_p(\beta_\theta(t_{p+2}) \circ \eta_{p-1}^-(t_1, \dots, t_{p-1}; t_{p+2})),
\end{aligned}$$

where $\eta_{p+1}^+(\cdot, r)$, $\eta_p^0(\cdot, r)$ and $\eta_{p-1}^-(\cdot, r)$ are the functions defined in the lemma with $f_{p+1}(\cdot, r)$ and β_r respectively in place of $f_p(\cdot)$ and β_θ .

From this, writing the kernels of orders p , $p+1$ and $p+2$, it may be concluded that the multiple integral kernels have the desired form. Note also that we have proved that $I_1(\beta_\theta)_{p+1}(f_{p+1})$ is in the domain of ∇^+ . □

Let Σ_p denote the set of all permutations of $\{1, \dots, p\}$. In this corollary we replace f_p by the product $h_1 \circ \dots \circ h_p$ to give a partial symmetrization of the kernels obtained in the above proposition.

Corollary 1 *If $f_p = h_1 \circ \dots \circ h_p$ where $h_i \in L^2(\mathbb{R}_+)$ for any $i = 1, \dots, p$ we have the following expression:*

$$I_1(\beta_\theta)I_p(h_1 \circ \dots \circ h_p) = I_{p+1}(g_{p+1}^+(\cdot; \theta)) + I_p(g_p^\circ(\cdot; \theta)) + I_{p-1}(g_{p-1}^-(\cdot; \theta))$$

where the symmetrizations of g_{p+1}^+ , g_p° and g_{p-1}^- are given by

$$\tilde{g}_{p+1}^+(t_1, \dots, t_{p+1}; \theta) = \sum_{\sigma \in \Sigma_p} \sum_{l=0}^p \frac{1}{l!} \beta_{t_{l+2}}(t_{l+1}) \prod_{i=1}^l h_{\sigma(i)}(t_i) \prod_{j=3}^{p+2-l} (\beta_{t_{l+j}} h_{\sigma(l+j-2)})(t_{l+j-1})$$

with $t_{p+2} = \theta$,

$$\begin{aligned}
&\tilde{g}_p^\circ(t_1, \dots, t_p; \theta) \\
&= \sum_{\sigma \in \Sigma_p} \sum_{l=0}^{p-1} \frac{1}{l!} (\alpha \beta_{t_{l+2}} h_{\sigma(l+1)})(t_{l+1}) \prod_{i=1}^l h_{\sigma(i)}(t_i) \prod_{j=3}^{p+1-l} (\beta_{t_{l+j}} h_{\sigma(l+j-1)})(t_{l+j-1})
\end{aligned}$$

with $t_{p+1} = \theta$, and

$$\tilde{g}_{p-1}^-(t_1, \dots, t_{p-1}; \theta) = \sum_{\sigma \in \Sigma_p} \sum_{l=0}^{p-1} \frac{1}{l!} \langle \beta_{t_{l+1}}, h_{\sigma(l+1)} \rangle \prod_{i=1}^l h_{\sigma(i)}(t_i) \prod_{j=2}^{p-l} (\beta_{t_{l+j}} h_{\sigma(l+j)})(t_{l+j-1})$$

with $t_p = \theta$.

We are now ready to state the chaotic Kabanov formula.

Theorem 1 *Let g be a bounded function with compact support on \mathbb{R}_+ , and*

$$\gamma(t_1, \dots, t_k) = \beta_{t_2}(t_1) \cdots \beta_{t_k}(t_{k-1}) g(t_k), \quad t_1, \dots, t_k \in \mathbb{R}_+, \quad k > 1,$$

and $\gamma(t) = g(t)$, $t \in \mathbb{R}_+$, for $k = 1$. We have

$$I_1(g)I_p(f_p) = I_{p+1}(g_{p+1}^+) + I_p(g_p^\circ) + I_{p-1}(g_{p-1}^-) \quad (18)$$

where the kernels g_{p+1}^+ , g_p° and g_{p-1}^- are

$$g_{p+1}^+(t_1, \dots, t_{p+1}) = \sum_{j=1}^{p+1} \frac{p!}{(j-1)!} \gamma(t_j, \dots, t_{p+1}) f_p(t_1, \dots, \hat{t}_j, \dots, t_{p+1}),$$

where \hat{t}_j means that t_j is omitted in the arguments of f_p ,

$$g_p^\circ(t_1, \dots, t_p) = f_p(t_1, \dots, t_p) \sum_{j=1}^p \frac{p!}{(j-1)!} \gamma(t_j, \dots, t_p) \alpha(t_j),$$

and

$$\begin{aligned} g_{p-1}^-(t_1, \dots, t_{p-1}) &= \sum_{j=1}^{p-1} \frac{p!}{(j-1)!} \gamma(t_j, \dots, t_{p-1}) \langle f_p(t_1, \dots, \dots, t_{p-1}, \cdot), \beta_{t_j}(\cdot) \rangle \\ &\quad + p \langle f_p(t_1, \dots, t_{p-1}, \cdot), g(\cdot) \rangle. \end{aligned}$$

Proof. By Proposition 2, we have

$$\begin{aligned} I_1(g)I_p(f_p) &= I_{p+1}(g \otimes f_p) + pI_{p-1}(\langle f_p(*, \cdot), g(\cdot) \rangle) \\ &\quad + p\nabla^+(g(\cdot)\alpha(\cdot)I_{p-1}(f_p(*, \cdot))) + p\nabla^+(g(\cdot)I_1(\beta)I_{p-1}(f_p(*, \cdot))), \end{aligned}$$

and applying Lemma 3 we obtain the decomposition (18).

□

Proposition 4 *Assume now that the underlying martingale is the standard Azéma martingale, i.e. $\alpha = 0$ and $\beta_\theta = -1_{[0,\theta]}$, $\theta \in \mathbb{R}_+$. Then*

$$I_1(1_{[0,\theta]})I_p(1_{[0,\theta]}^{\otimes p}) = \frac{1}{p+1}I_{p+1}(1_{[0,\theta]}^{\otimes(p+1)}) - pI_{p-1}((t_{(p-1)} - \theta)1_{[0,\theta]}^{\otimes(p-1)}(t_1, \dots, t_{p-1})),$$

with $t_{(0)} = 0$ and $t_{(p-1)} = t_1 \vee \dots \vee t_{p-1}$, $p > 1$.

Proof: We let $f_p(t_1, \dots, t_p) = \prod_{i=1}^p 1_{[0,\theta]}(t_i)$ and $g = 1_{[0,\theta]}$. By Th. 1 we have

$$\begin{aligned} g_{p+1}^+(t_1, \dots, t_{p+1}; \theta) &= \sum_{j=1}^{p+1} \frac{p!}{(j-1)!} (-1)^{p-j-1} 1_{\{t_j \leq \dots \leq t_{p+1} \leq \theta\}} \prod_{i=1}^{j-1} 1_{[0,\theta]}(t_i), \\ g_p^\circ(t_1, \dots, t_p; \theta) &= 0, \\ g_{p-1}^-(t_1, \dots, t_{p-1}; \theta) &= \sum_{j=1}^p (-1)^{p-j} \frac{p!}{(j-1)!} t_j 1_{\{t_j \leq \dots \leq t_{p-1} \leq \theta\}} \prod_{i=1}^{j-1} 1_{[0,\theta]}(t_i), \end{aligned}$$

with $t_p = \theta$.

Symmetrizing and doing some straightforward computations we obtain the desired result. □

5 Consequences of the chaotic Kabanov formula

In this section we use the chaotic Kabanov formula to prove that the Wiener and Poisson processes are the only normal martingales to possess certain properties relative to polynomials and Wick product.

5.1 Existence of a family of orthogonal polynomials associated to the martingale

Let X be a normal martingale. Following Meyer (1976) we define by induction the martingales

$$P_t^{(0)} = 1, \quad P_t^{(1)} = X_t, \dots, P_t^{(n)} = \int_{[0,t]} P_{s^-}^{(n-1)} dX_s, \quad t \in \mathbb{R}_+.$$

Note that $P_t^{(n)} = \frac{1}{n!} I_n(1_{[0,t]}^{\circ n})$ for all n . It is well known that in the Wiener case we have the relation $P_t^{(n)} = H_n(X_t, t)$, where $H_n(x, y) = y^{n/2} h_n\left(\frac{x}{\sqrt{y}}\right)$, and h_n is the Hermite polynomial of degree n .

Also in the Poisson case, we have $P_t^{(n)} = C_n(X_t, t)$, where $\{C_n(x, y)\}$ are the Charlier polynomials (Meyer, 1976), (Surgailis, 1984).

This situation motivates the following definition

Definition 3 *We will say that a normal martingale X has an associated family of polynomials $\{Q_n(x, y)\}$, where Q_n is a polynomial of degree n in x if*

$$P_t^{(n)} = Q_n(X_t, t), \text{ for all } n.$$

Now the natural problem is to characterize the normal martingales that have an associated family, and the answer is the next Theorem.

Theorem 2 *Let X be a normal martingale in $L^4(\Omega)$. Then X has an associated family of polynomials if and only if ϕ_s is a deterministic constant process, i.e. X is a Poisson or a Wiener process.*

Proof. The if part is straightforward. If $\phi_s = 0$ we are in the Wiener case, and if $\phi_s = c \neq 0$, we are in the Poisson case with jumps of height c .

The proof of the only if part is a consequence of the Kabanov formula. Assume that the chaotic decomposition of the process ϕ is

$$\phi_t = \sum_{n=0}^{\infty} I_n(f_n(\cdot, t)),$$

and X has an associated family of polynomials. Then by the Kabanov formula

$$\begin{aligned} P_t^{(1)} P_t^{(1)} &= X_t^2 = I_1(1_{[0,t]}) I_1(1_{[0,t]}) = I_2(1_{[0,t]}^{\circ 2}) + t + \nabla^+(1_{[0,t]}(x) \phi(x)) \\ &= I_2(1_{[0,t]}^{\circ 2}) + t + \sum_{i=1}^{\infty} I_i(1_{[0,t]}(x) f_{i-1}(\cdot, x)). \end{aligned} \quad (19)$$

For each fixed t , X_t^2 will be a linear combination of the three polynomials $Q_0(X_t, t)$, $Q_1(X_t, t)$ and $Q_2(X_t, t)$, which are respectively 1, $I_1(1_{[0,t]})$ and $\frac{1}{2} I_2(1_{[0,t]}^{\circ 2})$. So

$$X_t^2 = a_2(t)I_2(1_{[0,t]}^{\circ 2}) + a_1(t)I_1(1_{[0,t]}) + a_0(t). \quad (20)$$

Then the chaotic decomposition of X_t^2 has only terms until the second chaos, and from (19) we get that $f_n = 0$ for $n \geq 2$. Therefore

$$\phi_t = f_0(t) + I_1(f_1(\cdot, t)).$$

Identifying the kernels of (19) and (20), we have for all t

$$1_{[0,t]}(x)f_0(x) = a_1(t)1_{[0,t]}(x),$$

for the first chaos,

$$1_{[0,t]}^{\circ 2}(u, v) + \frac{1}{2} (1_{[0,t]}(u)f_1(v, u) + 1_{[0,t]}(v)f_1(u, v)) = a_2(t)1_{[0,t]}^{\circ 2}(u, v), \quad u, v \in \mathbb{R}_+,$$

for the second chaos. So $f_0(u) = c_0$, $u \in \mathbb{R}_+$, and $f_1(\cdot, x) = c_1 1_{[0,x]}(\cdot)$.

Then

$$\phi_t = c_0 + c_1 I_1(1_{[0,t]}), \quad t \in \mathbb{R}_+. \quad (21)$$

We apply now Proposition 2 (or the Kabanov formula) to the product $I_1(1_{[0,t]})I_2(1_{[0,t]}^{\circ 2})$ with a ϕ given by (21), and we get

$$\begin{aligned} I_1(1_{[0,t]})I_2(1_{[0,t]}^{\circ 2}) &= I_3(1_{[0,t]}^{\circ 3}) + 2tI_1(1_{[0,t]}) \\ &\quad + 2c_0I_2(1_{[0,t]}^{\circ 2}) + \nabla^+(2c_1 1_{[0,t]}(x_1)I_1(1_{[0,x_1]})I_1(1_{[0,t]})). \end{aligned}$$

Hence

$$\begin{aligned} I_1(1_{[0,t]})I_2(1_{[0,t]}^{\circ 2}) &= aI_3(1_{[0,t]}^{\circ 3}) + bI_2(1_{[0,t]}^{\circ 2}) \\ &\quad + cI_1 \left(t1_{[0,t]}(x_1) + \int_0^t c_1 1_{[0,x_1]}(u) du 1_{[0,t]}(x_1) \right), \end{aligned} \quad (22)$$

where $a, b, c \in \mathbb{R}$.

As

$$I_1(1_{[0,t]})I_2(1_{[0,t]}^{\circ 2}) = X_t Q_2(X_t, t),$$

the product is a polynomial of degree 3 in X_t , and can be expressed as a linear combination of Q_0 , Q_1 , Q_2 and Q_3 . So

$$\begin{aligned} I_1(1_{[0,t]})I_2(1_{[0,t]}^{\circ 2}) &= b_3(t)Q_3(X_t, t) + \dots + b_0(t) \\ &= \frac{1}{3!}b_3(t)I_3(1_{[0,t]}^{\circ 3}) + \frac{1}{2!}b_2(t)I_2(1_{[0,t]}^{\circ 2}) + b_1(t)I_1(1_{[0,t]}) + b_0(t), \end{aligned} \tag{23}$$

and as before identifying the kernels of (22) and (23), and focusing our attention to the chaos of order 1, we have

$$\begin{aligned} a_1(t)1_{[0,t]}(x_1) &= c \left(t1_{[0,t]}(x_1) + c_1 \int_0^t 1_{[0,x_1]}(u) du 1_{[0,t]}(x_1) \right) \\ &= c \left(t1_{[0,t]}(x_1) + c_1 x_1 1_{[0,t]}(x_1) \right), \end{aligned}$$

therefore c_1 has to be zero, and we get that the process ϕ is a constant. □

Remark 3 *Since for the standard Azéma martingale we have $\phi_s = -I_1(1_{[0,s]})$, the theorem implies that the Azéma martingale does not have an associated family of polynomials.*

Remark 4 *In particular we have proved that any normal martingale in L^4 with an associated family of polynomials is a process with independent increments. In fact, a normal martingale in L^4 has independent increments if and only if ϕ is deterministic (Utzet, 1992), (Emery, 1990).*

Remark 5 *Yor (1992, chapter 15) introduces a family of polynomials related to the Azéma martingales in a different context. These polynomials give the conditional expectation of powers for Azéma martingales with respect to the σ -algebras of the strict past.*

5.2 Projection property for the Wick product

Let X be a normal martingale in $L^4(\Omega)$. The Wick product $I_1(f_1) : I_n(g_n)$ of $I_1(f_1)$ and $I_n(g_n)$ is defined by

$$I_1(f_1) : I_n(g_n) = I_{n+1}(f_1 \circ g_n).$$

This motivates the following definition

Definition 4 We say that X has the Wick projection property if the Wick product $I_1(f_1) : I_n(g_n)$ is the projection of the product $I_1(f_1)I_n(g_n)$ over the chaos of order $n + 1$. Here the product $I_1(f_1)I_n(g_n)$ is assumed to belong to $L^2(\Omega)$ and the kernels f_1 and g_n are bounded and with compact support.

Now the natural problem is to determine the martingales that have this property.

Theorem 3 Assume that ϕ has a chaotic decomposition with no terms in chaos greater than 1. Then X has the Wick projection property if and only if ϕ is deterministic.

Proof. Proposition 2 says that

$$\begin{aligned} I_1(f_1)I_n(g_n) &= I_{n+1}(f_1 \circ g_n) + nI_{n-1} \left(\int_0^\infty f_1(s)g_n(\cdot, s) ds \right) \\ &\quad + \nabla^+(f_1(u)\phi(u)nI_{n-1}(g_n(\cdot, u))), \end{aligned}$$

and it is straightforward to see that the last addend has null projection over the chaos of order greater than n for all f_1 and g_n if and only if ϕ is deterministic. \square

Remark 6 If in the definition of the Wick projection property we impose $I_1(f_1) : I_n(g_n)$ to be the projection of the product $I_1(f_1)I_n(g_n)$ over the sum of chaos of order strictly greater than n , instead of over the class of order n , then the above result can be extended to normal martingales X with $\phi_t = \sum_i I_i(f_i(\cdot, t))$ for which Proposition 2 is valid.

6 Derivation rule of products

It is well-known that on the Wiener space, i.e. for $\phi = 0$, ∇^- is identified to a derivation operator. As noticed by Ma, Protter and San Martin (1998), ∇^- can not act in the same way as a Sobolev derivative, i.e. $\nabla^- f(M_t) = 1_{[0,t]}f'(M_t)$, unless $\phi = 0$. In this section we study the product rule for ∇^- and in particular we further show that ∇^- can be a derivation operator only for $\phi = 0$.

For deterministic ϕ it was noticed by Privault (1996) that the product formula becomes

$$\nabla_s^-(FG) = F\nabla_s^-G + G\nabla_s^-F + \phi_s\nabla_s^-F\nabla_s^-G, \quad s \in \mathbb{R}_+.$$

This rule does not extend to random ϕ , but (8) gives by duality the following information.

Proposition 5 For $F, G \in \mathcal{S}$,

$$E[\nabla_s^-(FG) \mid \mathcal{F}_s] = E[F\nabla_s^-G + G\nabla_s^-F + \phi_s\nabla_s^-F\nabla_s^-G \mid \mathcal{F}_s], \quad ds - a.e.$$

Proof. We write (8) for $u \in \mathcal{U}$ adapted, and apply the duality between ∇^- and ∇^+ :

$$\begin{aligned} E[(u, \nabla^-(FG))] &= E[\nabla^+(u)FG] \\ &= E[G(\nabla^+(u)F) + (u, \nabla^-F)_{L^2(\mathbb{R}_+)} + \nabla^+(u\phi\nabla^-F)] \\ &= E[(u, F\nabla^-G)_{L^2(\mathbb{R}_+)} + (u, G\nabla^-F)_{L^2(\mathbb{R}_+)} + (u, \phi\nabla^-F\nabla^-G)_{L^2(\mathbb{R}_+)}]. \quad \square \end{aligned}$$

From the above it follows that if ∇^- is a derivation then $E[\phi_s\nabla_s^-F\nabla_s^-G] = 0$, $s \in \mathbb{R}_+$, $F, G \in \mathcal{S}$, hence $\phi = 0$.

As a consequence of Proposition 5, we have

$$\nabla_s^-(FG) = F\nabla_s^-G + G\nabla_s^-F + \phi_s\nabla_s^-F\nabla_s^-G + \mathcal{A}_s(F, G), \quad s \in \mathbb{R}_+,$$

where $\mathcal{A}(F, G)$ is a process with zero adapted projection.

Our aim in the sequel is to gain more information on this process in the case of the Azéma martingales. We use the notation $\mathcal{A}_s(F, t) = \mathcal{A}_s(F, I_1(\beta_t))$, $t \in \mathbb{R}_+$.

Proposition 6 Assume that ϕ is given as $\phi_t = \alpha_t + I_1(\beta_t)$, $t \in \mathbb{R}_+$, where α is locally bounded and β_t is bounded, $t \in \mathbb{R}_+$. Then

$$\nabla_s^-(I_1(\beta_t)F) = F\nabla_s^-I_1(\beta_t) + I_1(\beta_t)\nabla_s^-F + \phi_s\nabla_s^-I_1(\beta_t)\nabla_s^-G + \mathcal{A}_s(F, t), \quad s \in \mathbb{R}_+,$$

where $E[\mathcal{A}_s(F, t) \mid \mathcal{F}_s] = 0$, $s \in \mathbb{R}_+$, and $\mathcal{A}_s(F, t)$ is given by the relation

$$\mathcal{A}_s(F, t) = \nabla^+(\beta_t(\cdot))\beta(s)\nabla^-F + \phi_s\beta_t(\cdot)\beta(s)\nabla^-\nabla_s^-F + \beta_t\mathcal{A}_s(\nabla^-F, \cdot). \quad (24)$$

Proof. We use the relation $\nabla_s^- \nabla^+(u) = \nabla^+(\nabla_s^- u) + u_s$:

$$\begin{aligned}
& \nabla_s^-(I_1(\beta_t)F) \\
&= \nabla_s^-(\nabla^+(\beta_t F) + (\beta_t, \nabla^- F)_{L^2(\mathbf{R}_+)} + \nabla^+(\beta_t \phi \nabla^- F)) \\
&= \nabla^+(\beta_t \nabla_s^- F) + \beta_t(s)F + (\beta_t, \nabla^- \nabla_s^- F)_{L^2(\mathbf{R}_+)} + \nabla^+(\beta_t \alpha \nabla^- \nabla_s^- F) \\
&\quad + \beta_t(s) \alpha_s \nabla_s^- F + \nabla^+(\beta_t(\cdot) \nabla_s^-(I_1(\beta_s) \nabla^- F)) + \beta_t(s) I_1(\beta_s) \nabla_s^- F \\
&= \nabla^+(\beta_t \nabla_s^- F) + (\beta_t, \nabla^- \nabla_s^- F)_{L^2(\mathbf{R}_+)} + \nabla^+(\beta_t \alpha \nabla^- \nabla_s^- F) \\
&\quad + F \nabla_s^- I_1(\beta_t) + \beta_t(s) \alpha_s \nabla_s^- F + \beta_t(s) I_1(\beta_s) \nabla_s^- F \\
&\quad + \nabla^+(\beta_t(\cdot) (\beta_s(s) \nabla^- F + I_1(\beta_s) \nabla^- \nabla_s^- F + \phi_s \nabla_s^- I_1(\beta_s) \nabla^- \nabla_s^- F + \mathcal{A}(\nabla^- F, \cdot))) \\
&= \nabla^+(\beta_t \nabla_s^- F) + (\beta_t, \nabla^- \nabla_s^- F)_{L^2(\mathbf{R}_+)} + \nabla^+(\beta_t \phi \nabla^- \nabla_s^- F) \\
&\quad + F \nabla_s^- I_1(\beta_t) + \beta_t(s) \alpha_s \nabla_s^- F + \beta_t(s) I_1(\beta_s) \nabla_s^- F \\
&\quad + \nabla^+(\beta_t(\cdot) \beta_s(s) \nabla^- F + \beta_t(\cdot) \phi_s \nabla_s^- I_1(\beta_s) \nabla^- \nabla_s^- F + \beta_t(\cdot) \mathcal{A}_s(\nabla^- F, \cdot)) \\
&= I_1(\beta_t) \nabla_s^- F + F \nabla_s^- I_1(\beta_t) + \phi_s \nabla_s^- I_1(\beta_t) \nabla_s^- F \\
&\quad + \nabla^+(\beta_t(\cdot) \beta_s(s) \nabla^- F + \beta_t(\cdot) \phi_s \beta_s(s) \nabla^- \nabla_s^- F + \beta_t(\cdot) \mathcal{A}_s(\nabla^- F, \cdot)),
\end{aligned}$$

and we obtain (24). □

As a consequence of this proposition, for $F = I_n(f^{on})$ the remaining process $\mathcal{A}(F, t)$ can be explicitly determined by induction from

$$\begin{aligned}
& \mathcal{A}_s(I_n(f^{on}), t) \\
&= n I_n((\beta_t(\cdot) \beta_s(s) f) \circ f^{\circ(n-1)}) + n(n-1) f_s \nabla^+(\phi_s \beta_t(\cdot) \beta_s(s) f I_{n-2}(f^{\circ(n-2)})) \\
&\quad + n \nabla^+(f \beta_t(\cdot) \mathcal{A}_s(I_{n-1}(f^{\circ(n-1)})), t).
\end{aligned}$$

Note that $F \mapsto \mathcal{A}(F, t)$ is linear and that the chaotic expansion of $\phi_s I_{n-2}(f^{\circ(n-2)})$ can be explicitly computed from the Kabanov formula. As an application we compute $\mathcal{A}_s(I_n(f^{on}), t)$ for $n = 0, 1, 2$. We have $\mathcal{A}_s(1, t) = 0$, $\mathcal{A}_s(I_1(f^{\circ 1}), t) = I_1(\beta_t(\cdot) \beta_s(s) f)$, and

$$\begin{aligned}
\mathcal{A}_s(I_2(f^{\circ 2}), t) &= I_2(\beta_t(\cdot) \beta_s(s) f \circ f) + 2 f_s \nabla^+((\alpha_s + I_1(\beta_s)) \beta_t(\cdot) \beta_s(s) f) \\
&\quad + 2 \nabla^+(f \beta_t(\cdot) I_1(\beta_s) \beta_s(s) f) \\
&= I_2(\beta_t(\cdot) \beta_s(s) f \circ f) + 2 f_s \alpha_s I_1(\beta_t(\cdot) \beta_s(s) f) + 2 f_s I_2(\beta_t(\cdot) \beta_s(s) f \circ \beta_s) \\
&\quad + 2 I_2(f \beta_t(\cdot) \circ \beta_s) \beta_s(s) f.
\end{aligned}$$

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