The Wiener chaos expansion for the Cox-Ingersoll-Ross model

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1. Introduction

We take the CIR model as given by

$$dr_t = a(b - r_t)dt + c\sqrt{r_t}d\tilde{W}_t, \qquad \|\lambda_t\|^2 = \bar{\lambda}^2 r_t \qquad (1)$$

for some positive constants $a,b,c,\bar{\lambda}$ with $4ab>c^2$, where \tilde{W}_t is a standard one dimensional Brownian motion on a probability space (Ω,\mathcal{F},P) equipped with a Brownian filtration $(\mathcal{F}_t)_{0\leq t\leq\infty}$.

• We seek for the chaotic representation of the underlying random variable X_{∞} in the Hughston/Rafailidis framework.

2. Positive Interest Rates

Let P_{tT} , $0 \le t \le T$ denote the price at time t for a zero coupon bond which pays one unit of currency at its maturity T. Clearly $P_{tt} = 1$ for all $0 \le t < \infty$ and furthermore, positivity of the interest rate is equivalent to having

$$P_{ts} \le P_{tu}, \tag{2}$$

for all $0 \le t \le u \le s$.

2.1 The Flesaker–Hughston approach

FH introduced zero coupon bond prices in the form

$$P_{tT} = \frac{\int_{T}^{\infty} h_s M_{ts} ds}{\int_{t}^{\infty} h_s M_{ts} ds}, \quad \text{for } 0 \le t \le T < \infty, \tag{3}$$

where M_{ts} is a family of strictly positive continuous martingales satisfying $M_{0s}=1$ and $h_T=-\frac{\partial P_{0T}}{\partial T}$ is a positive deterministic function obtained from the initial term structure. Then the positivity condition (2), as well as $P_{tt}=1$, holds for all $0 \le t \le T < \infty$.

For concrete examples, they introduce the process

$$V_t = \int_t^\infty h_s M_{ts} ds, \tag{4}$$

which is easily seen to be a strictly positive supermartingale.

Due to the martingale property for M_{ts} , the bond prices can be rewritten as

$$P_{tT} = \frac{E_t[V_T]}{V_t}. (5)$$

2.2 State price density and the potential approach

One can elevate equation (5) to the starting point of the modelling and concentrate on the positive adapted continuous process V_t , called the *state price density*.

Positivity of the interest rates is then equivalent to V_t being a supermartingale. In order to match the initial term structure, it needs to be chosen so that $E[V_T] = P_{0T}$. If we further impose that $P_{0T} \to 0$ as $T \to \infty$, then V_t satisfies all the properties of what is known in probability theory as a *potential* (namely, a positive supermartingale with expected value going to zero at infinity).

It follows from the Doob–Meyer decomposition that any continuous potential satisfying

$$E\left(\sup_{0\le t\le\infty} V_t^2\right) < \infty \tag{6}$$

can be written as

$$V_t = E_t[A_\infty] - A_t, \tag{7}$$

for a unique (up to indistinguishibility) adapted continuous increasing process A_t with $E(A_\infty^2)<\infty$, satisfying the constraint that

$$E\left[\frac{\partial A_T}{\partial T}\right] = -\frac{\partial P_{0T}}{\partial T}.$$
 (8)

2.3 Related quantities

Given a strictly positive supermartingale V_t , there exist a unique strictly positive (local) martingale Λ_t such that the process

$$B_t = \Lambda_t / V_t \tag{9}$$

is strictly increasing. We identify B_t with a riskless money market account initialized at $B_0=\mathbf{1}$ and write it as

$$B_t = \exp(\int_0^t r_s ds), \tag{10}$$

for an adapted process $r_s > 0$, the short rate process.

The market price of risk then arises as the adapted vector valued process λ_t such that

$$d\Lambda_t = -\lambda_t \Lambda_t dW_t. \tag{11}$$

It is immediate to see that

$$dV_t = -r_t V_t dt - \lambda_t V_t dW_t. \tag{12}$$

so that the specification of the process V_t is enough to produce both the short rate r_t and the market price of risk λ_t .

2.4 The Chaotic Approach

Assume that the state price density V_t is a potential satisfying

$$E\left[\int_0^\infty r_s V_s ds\right] < \infty \tag{13}$$

By integrating (12) on the interval (t,T), taking conditional expectations at time t and the limit $T \to \infty$, one finds that

$$V_t = E_t \left[\int_t^\infty r_s V_s ds \right] \tag{14}$$

Now let σ_t be a vector valued process such that

$$\|\sigma_t\|^2 = r_t V_t,\tag{15}$$

and define the square integrable random variable

$$X_{\infty} = \int_{0}^{\infty} \sigma_s dW_s. \tag{16}$$

It follows from the Ito isometry that

$$V_t = E_t[X_{\infty}^2] - E_t[X_{\infty}]^2, \tag{17}$$

which is called the conditional variance representation of the state price density V_t . A direct comparison between (14) and (4) gives that

$$h_s M_{ts} = E_t \left[\|\sigma_s\|^2 \right]. \tag{18}$$

Similarly, by comparing the conditional variance representation (17) with the decomposition (7), we see that

$$E_t[X_{\infty}^2] - X_t^2 = E_t[A_{\infty}] - A_t,$$

where $X_t = E_t[X_{\infty}]$. It follows from the uniqueness of the Doob-Meyer decomposition that

$$A_t = [X, X]_t,$$

that is, the quadratic variation of the process X_t .

2.5 Wiener chaos

Let W_t be an N-dimensional Brownian motion on the filtered probability space $(\Omega, \mathcal{F}, \{\mathcal{F}_t\}_{t \in \mathbb{R}_+}, P)$. To streamline the handling of vector indices by a compact notation

$$\tau = (s, \mu) \in \Delta \stackrel{\cdot}{=} \mathbb{R}_+ \times \{1, \dots, N\}$$

and express integrals as

$$\int_{\Delta} f(\tau)d\tau \doteq \sum_{\mu} \int_{0}^{\infty} f(s,\mu)ds$$

$$\int_{\Delta} f(\tau)dW_{\tau} \doteq \sum_{\mu} \int_{0}^{\infty} f(s,\mu)dW_{s}^{\mu} \tag{19}$$

For each $n \geq 0$, let

$$H_n(x) = (-1)^n e^{x^2/2} \frac{d^n}{dx^n} e^{-x^2/2}$$
 (20)

be the nth Hermite polynomial. For $h \in L^2(\Delta)$, let W(h) denote the Gaussian random variable $\int h(\tau)dW_{\tau}$. The spaces

$$\mathcal{H}_n \doteq \operatorname{span}\{H_n(W(h))|h \in L^2(\Delta)\}, \quad n \geq 1,$$

 $\mathcal{H}_0 \doteq \mathbb{C}$

form an orthogonal decomposition of the space $L^2(\Omega, \mathcal{F}_{\infty}, P)$ of square integrable random variables:

$$L^2(\Omega, \mathcal{F}_{\infty}, P) = \bigoplus_{n=0}^{\infty} \mathcal{H}_n$$

Each $\mathcal{H}^{(n)}$ can be understood completely via the *isometries*

$$f_n \mapsto J_n(f_n) = \int_{\Delta_n} f_n(\tau_1, \dots, \tau_n) dW_{\tau_1} \dots dW_{\tau_n}$$
 (21)

where $\Delta_n \doteq \{(\tau_1, ..., \tau_n) | \tau_i = (s_i, \mu_i) \in \Delta, 0 \le s_1 \le s_2 \le \cdots \le s_n < \infty \}$.

With these ingredients, one is then led to the result that any $X \in L^2(\Omega, \mathcal{F}_{\infty}, P)$ can be represented as a Wiener chaos expansion

$$X = \sum_{n=0}^{\infty} J_n(f_n) \tag{22}$$

where the deterministic functions $f_n \in L^2(\Delta_n)$ are uniquely determined by the random variable X.

A special example arises by noting that for $h \in L^2(\Delta)$

$$n!J_n(h^{\otimes n}) = ||h||^n H_n\left(\frac{W(h)}{||h||}\right)$$
 (23)

and furthermore

$$\exp\left[W(h) - \frac{1}{2} \int |h(\tau)|^2 d\tau\right] = \sum_{n=0}^{\infty} \frac{\|h\|^n}{n!} H_n\left(\frac{W(h)}{\|h\|}\right)$$
(24)

In the notation of quantum field theory, this example defines the Wick ordered exponential and Wick powers

$$: \exp[W(h)] : \doteq \exp\left[W(h) - \frac{1}{2} \int |h(\tau)|^2 d\tau\right]$$
$$: W(h)^n : \doteq n! J_n(h^{\otimes n}) \tag{25}$$

Theorem 1 For any random variable $X \in L^2(\Omega, \mathcal{F}_{\infty}, P)$, the generating functional $Z_X(h) : L^2(\Delta) \to \mathbb{C}$ defined by

$$Z_X(h) \doteq E\left[X \exp\left[W(h) - \frac{1}{2} \int |h(\tau)|^2 d\tau\right]\right]$$
 (26)

is an entire analytic functional of $h \in L^2(\Delta)$ and hence has an absolutely convergent expansion

$$Z_X(h) = \sum_{n \ge 0} F_X^{(n)}(h)$$
 (27)

where

$$F_X^{(n)}(h) = \int_{\Delta_n} f_X^{(n)}(\tau_1, \dots, \tau_n) h(\tau_1) \dots h(\tau_n) d\tau_1 \dots d\tau_n \qquad (28)$$

Here, $f_X^{(n)}(\tau_1, \ldots, \tau_n)$ is the nth Frechet derivative of Z_X at h=0. Finally, the Wiener–Itô chaos expansion of X is

$$X = \sum_{n>0} \int_{\Delta_n} f_X^{(n)}(\tau_1, \dots, \tau_n) dW_{\tau_1} \dots dW_{\tau_n}$$
 (29)

3. Squared Gaussian models

The CIR model with an integer constraint $N \doteq \frac{4ab}{c^2} \in \mathbb{N}_+ \setminus \{0,1\}$ lies in the class of so-called squared Gaussian models. By introducing an \mathbb{R}^N -valued Ornstein-Uhlenbeck process R_t , governed by the stochastic differential equation

$$dR_t = -\frac{a}{2}R_t dt + \frac{c}{2}dW_t \tag{30}$$

where W_t is N-dimensional Brownian motion, one verifies that the square $r_t = ||R_t||^2$ satisfies (1) where

$$\tilde{W}_t = \int_0^t ||R_t||^{-1} R_t \cdot dW_t$$

is itself a one-dimensional Brownian motion.

Definition 2 A pair (r_t, λ_t) of $(\Omega, \mathcal{F}, \mathcal{F}_t, P)$ processes is called an N-dimensional squared Gaussian model of interest rates $(N \geq 2)$ if there is an \mathbb{R}^N -valued Ornstein-Uhlenbeck process such that $r_t = ||R_t||^2$ and $\lambda_t = \bar{\lambda} R_t$. The process R_t satisfies

$$dR_t = \alpha(t)(\bar{R}(t) - R_t)dt + \gamma(t)dW_t, \quad R|_{t=0} = R_0$$
 (31)

where $\alpha, \gamma, \bar{\lambda}$ are symmetric matrix valued and \bar{R} vector valued deterministic measurable functions on \mathbb{R}_+ . W is standard N- dimensional Brownian motion. In addition we impose the boundedness condition:

• there is some constant M>0 such that $\alpha(t)\geq M$ and $|\bar{\lambda}(t)|^{-2}\geq M$ for all t

The exact solution of (31) is easily seen to be

$$R_t = \tilde{R}(t) + \int K(t, t_1) (\gamma dW)_{t_1}$$
 (32)

where

$$\tilde{R}(t) = K(t,0)R_0 + \int K(t,t_1)\alpha(t_1)\bar{R}(t_1)dt_1$$
 (33)

and $K(t,s), t \geq s$ is the matrix valued solution of

$$\begin{cases} dK(t,s)/dt = -\alpha(t)K(t,s) & 0 \le s \le t \\ K(t,t) = I & 0 \le t \end{cases}$$
(34)

which generates the Ornstein-Uhlenbeck semigroup.

By (12), the state price density process is

$$V_t = \exp\left[-\int_0^t \left(R_s^{\dagger} \left(1 + \frac{\bar{\lambda}^2}{2}\right) R_s ds + R_s^{\dagger} \bar{\lambda} dW_s\right)\right]$$
 (35)

We thus have a natural candidate for the random variable X_{∞} :

$$X_{\infty} = \int_0^{\infty} \sigma_t^{\dagger} \ dW_t \tag{36}$$

where σ_t is the \mathbb{R}^N -valued process

$$\sigma_t \doteq \exp\left[-\int_0^t \left(R_s^{\dagger} \left(\frac{1}{2} + \frac{\bar{\lambda}^2}{4}\right) R_s ds + \frac{1}{2} R_s^{\dagger} \bar{\lambda} dW_s\right)\right] R_t \tag{37}$$

is the natural solution of $\sigma_t^{\dagger} \sigma_t = r_t V_t$.

Proposition 3 $\lim_{T\to\infty} E[X_T^2] = 1$

4. Exponentiated second chaos

The chaos expansion we seek for the CIR model will be derived from a closed formula for expectations of e^{-Y} for elements

$$Y = A + \int_{\Delta} B^{\dagger}(\tau_1) dW_{\tau_1} + \int_{\Delta_2} C(\tau_1, \tau_2) dW_{\tau_1} dW_{\tau_2}$$
 (38)

in a certain subset $\mathcal{C}^+ \subset \mathcal{H}_{\leq 2} \doteq \mathcal{H}_0 \oplus \mathcal{H}_1 \oplus \mathcal{H}_2$.

If in (38) we define $C(\tau_1, \tau_2) = C(\tau_2, \tau_1)$ when $\tau_1 > \tau_2$, then C is the kernel of a symmetric integral operator on $L^2(\Delta)$:

$$[Cf](\tau) = \int_0^\infty C(\tau, \tau_1) f(\tau_1) d\tau_1 \tag{39}$$

Recall that integral Hilbert-Schmidt operators are finite norm operators under the norm:

$$||C||_{HS}^2 = \int_{\Delta^2} |C(\tau_1, \tau_2)|^2 d\tau_1 d\tau_2$$

We say that $Y \in \mathcal{H}_{\leq 2}$ is in \mathcal{C}^+ if C is the kernel of a symmetric Hilbert–Schmidt operator on $L^2(\Delta)$ such that (1+C) has nonnegative spectrum.

Proposition 4 Let $Y \in C^+$. Then

$$E[e^{-Y}] = [\det_2(1+C)]^{-1/2}$$

$$\exp\left[-A + \frac{1}{2} \int_{\Delta_2} B^{\dagger}(\tau_1)(1+C)^{-1}(\tau_1, \tau_2)B(\tau_2)d\tau_1 d\tau_2\right]$$

Remark 5 The Carleman–Fredholm determinant is defined as the extension of the formula

$$\det_2(1+C) = \det(1+C) \exp[-\text{Tr}(C)] \tag{40}$$

from finite rank operators to bounded Hilbert–Schmidt operators; the kernel $(1+C)^{-1}(\tau_1,\tau_2)$ is also the natural extension from the finite rank case.

Corollary 6 (Wick's theorem) The random variable $X = e^{-Y}$, for $Y = \int_{\Delta_2} C(\tau_1, \tau_2) dW_{\tau_1} dW_{\tau_2} \in \mathcal{C}^+$ has Wiener chaos coefficient functions

$$f_n(\tau_1, \dots, \tau_n) = \begin{cases} K \sum_{G \in \mathcal{G}_n} \prod_{g \in G} [C(1+C)^{-1}](\tau_{g_1}, \tau_{g_2}) & n \text{ even} \\ 0 & n \text{ odd} \end{cases}$$

where $K = [\det_2(1+C)]^{-1/2}$ and for n even, \mathcal{G}_n is the set of Feynman graphs on the n marked points $\{\tau_1, \ldots, \tau_n\}$. Each Feynman graph G is a disjoint union of unordered pairs $g = (\tau_{g_1}, \tau_{g_2})$ with $\bigcup_{g \in G} g = \{\tau_1, \ldots, \tau_n\}$.

Proof: The generating functional for $X = e^{-Y}$ is

$$Z_X(h) = E\left[X \exp\left(\int h(\tau)dW_{\tau} - \frac{1}{2}\int h(\tau)^2 d\tau\right)\right]$$

=
$$E\left[\exp\left(\int h(\tau)dW_{\tau} - \frac{1}{2}\int h(\tau)^2 d\tau - \int_{\Delta_2} C(\tau_1, \tau_2)dW_{\tau_1}dW_{\tau_2}\right)\right],$$

so we can use Proposition 4

$$Z_X(h) = \det_2(1+C)^{-1/2}$$

$$\exp\left[-\frac{1}{2}\int_{\Delta^2} h^{\dagger}(\tau_1)[\delta(\tau_1, \tau_2) - (1+C)^{-1}(\tau_1, \tau_2)]h(\tau_2)d\tau_1d\tau_2\right]$$

using the last part of Theorem 1, the result comes by evaluating the nth Fréchet derivative at h=0, or equivalently by expanding the exponential and symmetrizing over the points τ_1,\ldots,τ_n in the n/2th term.

5. The chaotic expansion for squared Gaussian models

We now derive the chaos expansion for the squared Gaussian model defined by (31). For simplicity, we assume from now on that $\bar{\lambda}=0$. In view of (36) it will be enough to find the chaos expansion for σ_T^μ , $T<\infty$. We start by finding its the generating functional $Z_{\sigma_T^\mu}$. Define the functional $Z(h,k)=E\left[e^{-Y_T}\right]$ with

$$Y_{T} = \frac{1}{2} \int_{0}^{T} R_{t}^{\dagger} R_{t} dt + - \int_{0}^{T} h^{\dagger}(t) dW_{t} + \frac{1}{2} \int_{0}^{T} \|h(t)\|^{2} dt - \int_{0}^{T} k^{\dagger}(t) R_{t} dt.$$
 (41)

Proposition 7 Z(h,k) is an entire analytic functional of $(h,k) \in L^2([0;T]) \otimes \mathbb{R}^{2N}$. Moreover

$$\lim_{t \to T^{-}} \frac{\delta Z(h,k)}{\delta k^{\mu}(t)} \Big|_{k=0} = Z_{\sigma_T^{\mu}}(h)$$
(42)

where $Z_{\sigma_T^{\mu}}(h)$ is defined by (26) with $X=\sigma_T^{\mu}, \mu=1,\ldots,N$.

We want to use Proposition 4 in order to compute Z(h,k). Substitution of (32) into the first and last terms of (41) leads to

$$\frac{1}{2} \int_{0}^{T} R_{t}^{\dagger} R_{t} dt = \frac{1}{2} \int_{0}^{T} \tilde{R}^{\dagger}(t) \tilde{R}(t) dt + \int_{0}^{T} \left(\int_{0}^{T} \tilde{R}_{s}^{\dagger} K_{T}(s, t) ds \right) \gamma(t) dW_{t}
+ \int_{\Delta_{2}} \gamma(t_{1}) \left(\int_{0}^{T} K_{T}^{\dagger}(s, t_{1}) K_{T}(s, t_{2}) ds \right) \gamma(t_{2}) dW_{t_{1}} dW_{t_{2}}
+ \frac{1}{2} \int_{0}^{T} \operatorname{tr} \left[\gamma(t) \left(\int_{0}^{T} K_{T}^{\dagger}(s, t) K_{T}(s, t) ds \right) \gamma(t) \right] dt,
\int_{0}^{T} k_{t}^{\dagger} R_{t} dt = \int_{0}^{T} k_{t}^{\dagger} \tilde{R}_{t} dt + \int_{0}^{T} \left(\int_{0}^{T} k_{s}^{\dagger} K_{T}(s, t) ds \right) \gamma(t) dW_{t},$$

where we define $K_T(t_1, t_2) = \mathbf{1}(t_1 \leq T)K(t_1, t_2)$ and use an operator multiplication notation.

Thus the exponent Y_T appearing in (41) has the form of (38) with

$$A_{T} = \frac{1}{2} \int_{0}^{T} \operatorname{tr} \left[\gamma(t) \left(\int_{0}^{T} K_{T}^{\dagger}(s, t) K_{T}(s, t) ds \right) \gamma(t) \right] dt$$

$$\frac{1}{2} \int_{0}^{T} \left[\tilde{R}_{t}^{\dagger} \tilde{R}_{t} + h_{t}^{\dagger} h_{t} - 2k_{t}^{\dagger} \tilde{R}_{t} \right] dt,$$

$$B_{T}(t) = -h_{t} - \gamma(t) \int_{0}^{T} K_{T}^{\dagger}(s, t) k_{s} ds + \gamma(t) \int_{0}^{T} K_{T}^{\dagger}(s, t) \tilde{R}_{s} ds,$$

$$C_{T}(t_{1}, t_{2}) = \gamma(t_{1}) \left(\int_{0}^{T} K_{T}^{\dagger}(s, t_{1}) K_{T}(s, t_{2}) ds \right) \gamma(t_{2}).$$

Now we note that $C_T(\cdot,\cdot)$ is the kernel of an integral operator which is manifestly positive and has Hilbert-Schmidt norm $\|C_T\|_{HS}^2 = \mathcal{O}(T)$. Therefore, we can use Proposition 4 for $E[e^{-Y_T}]$, leading to a general formula for the generating functional Z(h,k)

Differentiation once with respect to k then yields

$$Z_{\sigma_T}(h) = M_T \exp\left[-\frac{1}{2} \int_0^T \left(\tilde{R}_t^{\dagger} \tilde{R}_t + h_t^{\dagger} h_t\right) dt\right]$$

$$\left[-\tilde{R} + K_T \gamma (1 + C_T)^{-1} (h - \gamma K_T^{\dagger} \tilde{R})\right] (T)$$

$$\exp\left[\frac{1}{2} \int_{\Delta_2} \left(h^{\dagger} - \tilde{R}^{\dagger} K_T \gamma\right) (1 + C_T)^{-1} \left(h - \gamma K_T^{\dagger} \tilde{R}\right) dt_1 dt_2\right]$$

where

$$M_T = e^{-\frac{1}{2} \operatorname{tr} C_T} (\det_2(1 + C_T))^{-1/2} = (\det(1 + C_T))^{-1/2}$$
 (43)

For the CIR model with $r_0=0$ we have $\tilde{R}=0$, $\alpha(t)=a/2$ and $\gamma(t)=c/2$, so that $K_T(s,t)=e^{-a(s-t)/2}\mathbf{1}(t\leq s\leq T)$ and

$$C_T(t_1, t_2) = \frac{c^2}{4} \int_0^T K_T(s, t_1) K_T(s, t_2) ds = \frac{c^2}{4a} e^{\frac{a}{2}(t_1 + t_2)} (e^{-at_2} - e^{-aT}),$$
(44)

for $0 \le t_1 \le t_2 \le T$.

Moreover, the previous expression for $Z_{\sigma_T}(h)$ reduces to

$$Z_{\sigma_T}(h) = M_T \left[K_T \gamma (1 + C_T)^{-1} h \right] (T)$$

$$\exp \left[-\frac{1}{2} \int_0^T h_t^{\dagger} h_t dt + \frac{1}{2} \int_{\Delta_2} h_{t_1}^{\dagger} (1 + C_T)^{-1} (t_1, t_2) h_{t_2} dt_1 dt_2 \right].$$

Theorem 8 The nth term of the chaos expansion of σ_T for the CIR model with $\bar{\lambda}=0$ and initial condition $r_0=0$ is zero for n even. For n odd, the kernel of the expansion is the function $f_T^{(n)}(\cdot):\Delta_n\to\mathbb{R}$

$$f_T(t_1, \dots, t_n) = M_T \sum_{G \in \mathcal{G}_n^*} \prod_{g \in G} L(g)$$
(45)

where

$$L(g) = \begin{cases} [C_T(1+C_T)^{-1}](t_{g_1}, t_{g_2}) & t \notin g\\ (K_T\gamma(1+C_T)^{-1})(T, t_{g_2}) & t \in g \end{cases}$$
(46)

Here, \mathcal{G}_n^* is the set of Feynman graphs, each Feynman graph G being a partition of $\{t_1, \ldots, t_n, T\}$ into pairs $g = (t_{g_1}, t_{g_2})$.

6. Bond pricing formula

In this section we give a derivation of the price of a zero coupon bond in the CIR model. Recall from section 2 that these are given by

$$P_{tT} = E_t[V_t^{-1}V_T]. (47)$$

To mantain things as simple as possible we keep $\bar{\lambda}=0$ so $V_t=\exp[-\int_0^T r_s ds]$, or in terms of the squared Gaussian formulation,

$$V_t = \exp\left[-\int_0^t R_s^{\dagger} R_s ds\right]. \tag{48}$$

But as we have seen in the previous section, for $t \leq s \leq T$

$$R_s^{\mu} = K_T(s,t)R_t^{\mu} + \frac{c}{2} \int_t^s K_T(s,s_1) dW_{s_1}^{\mu}.$$
 (49)

Hence $-\log[V_t^{-1}V_T] = \sum_{\mu} \int_t^T (R_s^{\mu})^2 ds$ can be written as

$$\sum_{\mu} \left[\frac{1}{2} (R_t^{\mu})^2 C_T(t,t) + \frac{cR_t^{\mu}}{2} \int_t^T C_T(t,s) dW_s^{\mu} \right] + \int_t^T \int_t^{s_2} \frac{c^2}{4} C_T(s_1,s_2) dW_{s_1}^{\mu} dW_{s_2}^{\mu} \right] + N \int_t^T C_T(s,s) ds$$

where $C_T(s_1, s_2) = 2 \int_0^T K_T(s, s_1) K_T(s, s_2) ds$

Taking the conditional expectation of $V_t^{-1}V_T$ by use of Proposition 4 leads to the desired formula

$$P_{tT} = \left(\det(1 + \frac{c^2}{4}C_T) \right)^{-N/2} \tag{50}$$

$$\prod_{\mu} \exp\left[-\frac{1}{2}(R_t^{\mu})^2 \left(C_T (1 + \frac{c^2}{4}C_T)^{-1}\right) (t, t)\right]$$
 (51)

Thus P_{tT} has the exponential affine form $\exp[-\beta(t,T)r_t - \alpha(t,T)]$ with

$$\beta(t,T) = \frac{1}{2} [C_T (1 + \frac{c^2}{4} C_T)^{-1}](t,t)$$

$$\alpha(t,T) = \frac{N}{2} \log \left(\det(1 + \frac{c^2}{4} C_T) \right)$$
(52)

The known formula has the same form, with

$$\beta(t,T) = \frac{2(e^{\rho(T-t)} - 1)}{(\rho + a)(e^{\rho(T-t)} - 1) + 2\rho}, \quad \rho^2 = a^2 + 2c^2$$

$$\alpha(t,T) = -\frac{2ab}{\rho^2} \log \left[\frac{2\rho e^{(a+\rho)(T-t)/2}}{(\rho + a)(e^{\rho(T-t)} - 1) + 2\rho} \right]$$
(53)

which can be derived as solutions of the pair of Ricatti ordinary differential equations

$$\frac{\partial \beta}{\partial t} = \frac{c^2 \beta^2}{2} + a\beta - 1$$

$$\frac{\partial \alpha}{\partial t} = -ab\beta$$
(54)

7. Discussion

- The CIR model, at least in integer dimensions, can be viewed within the chaos framework of Hughston and Rafailidis as arising from a random variable X_{∞} derived from exponentiated second chaos random variables $e^{-Y}, Y \in \mathcal{C}^+$. Such exponentiated \mathcal{C}^+ variables form a rich and natural family which is likely to include many more candidates for applicable interest rate models.
- Although their analytic properties are complicated, there do exist approximation schemes which can in principle be the basis for numerical methods

• This family is invariant under conditional \mathcal{F}_t —expectations: $\log E_t[e^{-Y}] \in \mathcal{C}^+$ whenever $Y \in \mathcal{C}^+$.

• Illustration of the connection between methods developed for quantum field theory and the methods of Malliavin calculus.