

# Infinite dimensional Kolmogorov operators with time dependent drift coefficients

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Keywords: Kolmogorov operator,  $m$ -dissipativity, stochastic partial differential equation

The main object of this paper is the Kolmogorov operator in  $[0, T] \times H$ , where  $H$  is a separable Hilbert space with norm  $|\cdot|$  and inner product  $\langle \cdot, \cdot \rangle$  and  $T > 0$  is fixed, defined by

$$L_0 u = D_t u + N(t)u,$$

$$N(t)u(t, x) = \frac{1}{2} \operatorname{Tr} [CD_x^2 u(t, x)] + \langle x, A^* D_x u(t, x) \rangle + \langle F(t, x), D_x u(t, x) \rangle.$$

Here  $A: D(A) \subset H \rightarrow H$  is the infinitesimal generator of a  $C_0$  semigroup  $e^{tA}$  in  $H$ ,  $A^*$  is the adjoint of  $A$ ,  $C$  is a positive linear operator on  $H$  and the mapping  $F: D(F) \subset [0, T] \times H \rightarrow H$  is such that  $F(t, \cdot)$  is quasi dissipative for all  $t \in [0, T]$  (see the definitions below). By  $D_t$  and  $D_x$  we denote the derivatives in  $t$  and  $x$  respectively. In order to define the domain of  $L_0$  we introduce some functional spaces. We denote by  $\mathcal{E}_A(H)$  the linear span of all real and imaginary parts of functions  $e^{i\langle x, h \rangle}$  where  $h \in D(A^*)$ . For any  $\varphi \in C^1([0, T])$  such that  $\varphi(T) = 0$  and any  $h \in C^1([0, T]; D(A^*))$  let us consider the function

$$u_{\varphi, h}(t, x) = \varphi(t)e^{i\langle x, h(t) \rangle}, \quad t \in \mathbb{R}, x \in H,$$

and denote by  $\mathcal{E}_A([0, T] \times H)$  the linear span of all real and imaginary parts of such functions  $u_{\varphi, h}$ . The operator  $L_0$  is defined on the space  $D(L_0) := \mathcal{E}_A([0, T] \times H)$ . Let  $\mathcal{P}(H)$  be the set of all Borel probability measures on  $H$  and let  $L(H)$  be the space of bounded linear operators on  $H$ .

First we shall assume that  $F$  is sufficiently regular and prove that, for any  $\nu_0 \in \mathcal{P}(H)$ , there exists a unique family of probability measures  $(\nu_t)_{t \in [0, T]} \subset \mathcal{P}(H)$  such that

$$\frac{d}{dt} \int_H u(t, x) \nu_t(dx) = \int_H L_0 u(t, x) \nu_t(dx), \quad u \in D(L_0), t \in [0, T]. \quad (1)$$

This statement is an infinite dimensional analog of some results in [1]–[3]. Then we show that  $L_0$  is essentially  $m$ -dissipative in the space  $L^2([0, T] \times H; \nu)$  where  $\nu$  is the measure in  $[0, T] \times H$  defined by  $\nu(dt, dx) = \nu_t(dx)dt$ . In the case of irregular drifts we prove an analogous result under the assumption that there exists a suitable family of probability measures  $\nu_t$  (see Condition (H2)). Finally, we apply the obtained results to reaction-diffusion equations with time-dependent coefficients.

Let us list the assumptions on the linear operator  $A$  which we will assume throughout.

**Condition (H1):**

- (i) There is  $\omega > 0$  such that  $\langle Ax, x \rangle \leq -\omega|x|^2$ ,  $x \in D(A)$ ;
- (ii)  $C \in L(H)$  is symmetric, nonnegative and such that the linear operator

$$Q_t := \int_0^t e^{sA} C e^{sA^*} ds$$

is of trace class for all  $t > 0$ ;

- (iii)  $e^{tA}(H) \subset Q_t^{1/2}(H)$  for all  $t > 0$  and there is a bounded operator  $\Lambda_t$  such that  $Q_t^{-1/2} \Lambda_t = e^{tA}$  and

$$\gamma_\lambda := \int_0^{+\infty} e^{-\lambda t} \|\Lambda_t\| dt < +\infty.$$

Let us note that assumption (iii) implies that the Ornstein–Uhlenbeck operator associated to  $L_0$  (i.e., corresponding to  $F = 0$ ) is strong Feller. This assumption is not essential but it allows to simplify several proofs. The Ornstein–Uhlenbeck operator is defined by

$$U\varphi(x) = \frac{1}{2} \operatorname{Tr} [CD_x^2\varphi(x)] + \langle x, A^*D_x\varphi(x) \rangle, \varphi \in \mathcal{E}_A(H).$$

In addition, we introduce the operator

$$V_0u(t, x) = D_tu(t, x) + Uu(t, x), \quad u \in \mathcal{E}_A([0, T] \times H),$$

and its maximal extension  $V$ .

Assume first that, in addition to Condition (H1), the mapping  $F: [0, T] \times H \rightarrow H$  is continuous along with  $D_xF$  and there exists  $K > 0$  such that

$$|F(t, x) - F(t, y)| \leq K|x - y|, \quad x, y \in H, \quad t \in [0, T].$$

It is known (see, e.g., [4]) that, under this assumption, for any  $s \geq 0$ , there exists a unique mild solution  $X(\cdot, s, x)$  of the stochastic differential equation

$$dX = (AX + F(t, X))dt + \sqrt{C} dW(t), \quad X(s) = x \in H, \quad (2)$$

where  $W(t)$ ,  $t \geq 0$ , is a cylindrical Wiener process in  $H$  defined on a filtered probability space  $(\Omega, \mathcal{F}, \mathcal{F}_t, \mathbb{P})$ . A mild solution  $X(t, s, x)$  of (2) is an adapted stochastic process  $X \in C([s, T]; L^2(\Omega, \mathcal{F}, \mathbb{P}))$  such that

$$X(t, s, x) = e^{(t-s)A}x + \int_s^t e^{(t-r)A}F(r, X(r, s, x))dr + W_A(t, s), \quad t \geq s,$$

where  $W_A(t, s)$  is the *stochastic convolution*:

$$W_A(t, s) = \int_s^t e^{(t-r)A}\sqrt{C} dW(r), \quad t \geq s.$$

In view of Condition (H1)-(ii),  $W_A(t, s)$  is a Gaussian random element in  $H$  with mean 0 and covariance operator  $Q_{s,t}$  given by

$$Q_{s,t}x = \int_s^t e^{sA}C e^{sA^*} x ds, \quad t \geq s, \quad x \in H.$$

We define the transition evolution operator

$$P_{s,t}\varphi(x) = \mathbb{E}[\varphi(X(t, s, x))], \quad t \geq s, \quad \varphi \in C_b(H),$$

where  $C_u(H)$  is the Banach space of all uniformly continuous and bounded functions  $\varphi: H \rightarrow \mathbb{R}$  endowed with the usual supremum norm

$$\|\varphi\|_0 = \sup_{x \in H} |\varphi(x)|.$$

The spaces  $C_u^k(H)$ ,  $k \in \mathbb{N}$ , are defined similarly.

We denote by  $C_u^*(H)$  the topological dual of  $C_u(H)$ . For any  $0 \leq s < t \leq T$ , let  $P_{s,t}^*$  be the adjoint operator of  $P_{s,t}$ . It is easy to see that if  $\nu \in \mathcal{P}(H)$  we have  $P_{s,t}^*\nu \in \mathcal{P}(H)$  and

$$\int_H \varphi(x)(P_{s,t}^*\nu)(dx) = \int_H P_{s,t}\varphi(x)\nu(dx), \quad \forall \varphi \in C_u(H).$$

In the case under consideration we have the following existence and uniqueness results and moment estimates.

**Proposition 1.** *Let  $\nu_0 \in \mathcal{P}(H)$ . Then  $\nu_t = P_{0,t}^*\nu_0$  satisfies (1).*

**Proposition 2.** *Let  $(\zeta_t)_{t \in [0, T]}$  be a solution of (1) such that*

$$\sup_{t \in [0, T]} \int_H |x|^2 \zeta_t(dx) < +\infty.$$

*Then  $\zeta_t = P_{0,t}^*\nu_0$  for all  $t \in [0, T]$ .*

**Proposition 3.** *Let  $\nu_0 \in \mathcal{P}(H)$  and let  $\nu_t = P_{0,t}^* \nu_0$ ,  $t \in [0, T]$ . Then for any  $m \in \mathbb{N}$  there exists  $c_m > 0$  such that*

$$\int_H |x|^{2m} \nu_t(dx) \leq c_m \left( 1 + \int_H |x|^{2m} \nu(dx) \right), \quad t \geq 0.$$

We recall that a mapping  $T: D(T) \subset E \rightarrow E$ , where  $E$  is a Hilbert space, is called dissipative if  $\langle T(x) - T(y), x - y \rangle \leq 0$  for all  $x, y \in D(T)$ . If, in addition,  $(I - T)(D(T)) = E$ , then  $T$  is called  $m$ -dissipative. The mapping  $T$  is called quasi dissipative if  $T - KI$  is dissipative for some  $K > 0$ . If  $T - KI$  is  $m$ -dissipative, then  $T$  is called  $m$ -quasi dissipative.

It is readily verified that  $L_0$  is dissipative in  $L^2([0, T] \times H, \nu)$ , hence it is closable. Let us denote its closure by  $L_2$ .

**Theorem 1.** *The operator  $L_2$  is  $m$ -dissipative in the space  $L^2([0, T] \times H; \nu)$ .*

Let us turn to the case of irregular drifts. Suppose we are given a family  $\{F(t, \cdot)\}_{t \in [0, T]}$  of  $m$ -quasi dissipative mappings  $F(t, \cdot): D(F(t, \cdot)) \subset H \rightarrow H$ , where  $D(F(t, \cdot))$  are Borel sets in  $H$ . For simplicity we shall assume that these mappings are  $m$ -dissipative. We are concerned with the Kolmogorov operator

$$L_0 u(t, x) := D_t u(t, x) + U u(t, x) + \langle F_0(t, \cdot), D_x u(t, x) \rangle, \quad u \in D(L_0),$$

where  $D(L_0) = \mathcal{E}_A([0, T] \times H)$  and  $U$  is the Ornstein–Uhlenbeck operator. Our goal is to prove that the closure of  $L_0$  is  $m$ -dissipative in the space  $L^2([0, T] \times H, \nu)$ , where  $\nu(dt, dx) = \nu_t(dx) dt$  and  $(\nu_t)_{t \in [0, T]}$  is a given family of Borel measures on  $H$  such that

$$\frac{d}{dt} \int_H u(t, x) \nu_t(dx) = \int_H L_0 u(t, x) \nu_t(dx), \quad \forall u \in D(L_0). \quad (3)$$

In addition to Condition (H1) we shall assume also

**Condition (H2):**

(i) There is a family  $\{F(t, \cdot)\}_{t \in [0, T]}$  of  $m$ -quasi dissipative mappings in  $H$  such that  $0 \in D(F(t, \cdot))$  and  $F_0(t, 0) = 0$  for all  $t \in \mathbb{R}$ ;

(ii) there is a family  $(\nu_t)_{t \in [0, T]}$  of Borel probability measures on  $H$  such that

$$\int_H |x|^2 \nu_0(dx) < +\infty$$

and

$$\int_0^T \int_H (|x|^4 + |F_0(t, x)|^2 + |x|^4 |F_0(t, x)|^2) \nu_t(dx) dt < +\infty;$$

(iii) for all  $u \in D(L_0)$  we have  $L_0 u \in L^2([0, T] \times H, \nu)$  and (3) is fulfilled;

(iv)  $\nu_t(D(F(t, \cdot))) = 1$ ,  $\forall t \in [0, T]$ .

The problem of existence of measures  $\nu_t$  as above is studied in our separate paper.

Let us assume that Conditions (H1) and (H2) hold.

**Proposition 4.** *For all  $u \in D(L_0)$  we have*

$$\begin{aligned} \int_0^T \int_H L_0 u(t, x) u(t, x) \nu_t(dx) dt &= -\frac{1}{2} \int_0^T \int_H |C^{1/2} D_x u(t, x)|^2 \nu_t(dx) dt \\ &\quad - \int_H u^2(0, x) \nu_0(dx). \end{aligned}$$

*In particular,  $L_0$  is dissipative in  $L^2([0, T] \times H, \nu)$ .*

Since  $L_0$  is dissipative, it is closable in  $L^2([0, T] \times H, \nu)$ . We shall denote its closure with domain  $D(L_2)$  by  $L_2$ .

Now we can state the main result of this note.

**Theorem 2.** *Under Conditions (H1) and (H2), the operator  $L_2$  is  $m$ -dissipative in the space  $L^2([0, +\infty) \times H, \nu)$ .*

Now we apply our result to the reaction-diffusion equation. Let us consider a stochastic heat equation perturbed by a time dependent polynomial drift of odd degree  $d > 1$  of the form  $\lambda\xi - p(t, \xi)$ ,  $\xi \in \mathbb{R}$ ,  $t \in [0, T]$ , where  $\lambda \in \mathbb{R}$  is given and  $\partial_\xi p(t, \xi) \geq 0$  for all  $\xi \in \mathbb{R}$  and  $t \in [0, T]$ . We set  $H = L^2(\mathcal{O})$ , where  $\mathcal{O} = (0, 1)^n$ ,  $n \in \mathbb{N}$ , and denote by  $\partial\mathcal{O}$  the boundary of  $\mathcal{O}$ . Let us consider the following stochastic partial differential equation on  $\mathcal{O}$ :

$$dX(t, s, \xi) = [\Delta_\xi X(t, s, \xi) + \lambda X(t, s, \xi) - p(t, X(t, s, \xi))]dt + BdW(t, \xi), \quad t \geq s, \quad (4)$$

$$X(t, s, \xi) = 0, \quad t \geq s, \quad \xi \in \partial\mathcal{O}, \quad X(s, s, \xi) = x(\xi), \quad \xi \in \mathcal{O}, \quad x \in H,$$

where  $\Delta_\xi$  is the Laplace operator,  $B \in L(H)$ , and  $W$  is a cylindrical Wiener process in  $H$  defined on a probability space  $(\Omega, \mathcal{F}, \mathbb{P})$ . We choose  $W$  of the form

$$W(t, \xi) = \sum_{k=1}^{\infty} e_k(\xi) \beta_k(t), \quad \xi \in \mathcal{O}, \quad t \geq 0,$$

where  $(e_k)$  is a complete orthonormal system in  $H$  and  $(\beta_k)$  is a sequence of independent standard Brownian motions on a probability space  $(\Omega, \mathcal{F}, \mathbb{P})$ .

In order to write problem (4) as a stochastic differential equation in the Hilbert space  $H$  we denote by  $A$  the realization of the Laplace operator with Dirichlet boundary conditions, i.e.,

$$Ax = \Delta_\xi x, \quad x \in D(A), \quad D(A) = H^2(\mathcal{O}) \cap H_0^1(\mathcal{O}).$$

The operator  $A$  is self-adjoint and possesses a complete orthonormal system of eigenfunctions, namely

$$e_k(\xi) = (2/\pi)^{n/2} \sin(\pi k_1 \xi_1) \cdots (\sin \pi k_n \xi_n), \quad \xi = (\xi_1, \dots, \xi_n) \in \mathbb{R}^n,$$

where  $k = (k_1, \dots, k_n)$ ,  $k_i \in \mathbb{N}$ .

For any  $x \in H$  we set  $x_k = \langle x, e_k \rangle$ ,  $k \in \mathbb{N}^n$ . Notice that

$$Ae_k = -\pi^2 |k|^2 e_k, \quad k \in \mathbb{N}^n, \quad |k|^2 = k_1^2 + \cdots + k_n^2.$$

Therefore, we have

$$\|e^{tA}\| \leq e^{-\pi^2 t}, \quad t \geq 0.$$

Concerning the operator  $B$ , we shall assume for simplicity that  $B = (-A)^{-\gamma/2}$  with  $n/2 - 1 < \gamma < 1$  (which implies  $n < 4$ ). Now it is easy to check that Condition (H1) is fulfilled. In fact we have

$$\begin{aligned} Q_t x &= \int_0^t e^{sA} B B^* e^{sA^*} x ds = \int_0^t (-A)^{-\gamma} e^{2tA} x dt \\ &= (-A)^{-(1+\gamma)} (1 - e^{2tA}) x, \quad t \geq 0, \quad x \in H. \end{aligned}$$

Then

$$\text{Tr} [(-A)^{-(1+\gamma)}] = \pi^{-2(1+\gamma)} \sum_{k \in \mathbb{N}^n} |k|^{-2(1+\gamma)} < +\infty,$$

since  $\gamma > \frac{n}{2} - 1$ .

Now, setting  $X(t, s) = X(t, s, \cdot)$  and  $W(t) = W(t, \cdot)$ , we can write problem (4) as

$$dX(t, s) = [AX(t, s) + F(t, X(t, s))]dt + (-A)^{-\gamma/2} dW(t), \quad t \geq s, \quad X(s, s) = x, \quad (5)$$

where  $F$  is the mapping

$$F: D(F) = [0, T] \times L^{2d}(\mathcal{O}) \subset [0, T] \times H \rightarrow H, \quad x(\xi) \mapsto \lambda\xi - p(t, x(\xi)).$$

It is convenient, following [4], to introduce two different notions of solution of (5). For this purpose, for any  $s \in [0, T)$ , we consider the space

$$C_W([s, T]; L^2(\Omega, \mathcal{F}, \mathbb{P}; H)) := C_W([s, T]; H)$$

consisting of all continuous mappings  $F: [s, T] \rightarrow L^2(\Omega, \mathcal{F}, \mathbb{P}; H)$  adapted to  $W$ ; endowed with the norm

$$\|F\|_{C_W([s, T]; H)} = \left( \sup_{t \in [s, T]} \mathbb{E} (|F(t)|^2) \right)^{1/2}$$

$C_W([s, T]; H)$  is a Banach space.

**Definition 1.** (i) Let  $x \in L^{2d}(\mathcal{O})$ . We say that a mapping  $X(\cdot, s, x) \in C_W([s, T]; H)$  is a mild solution of problem (5) if  $X(t, s, x) \in L^{2d}(\mathcal{O})$  for all  $t \in [s, T]$  and the following integral equality holds:

$$X(t, s, x) = e^{(t-s)A}x + \int_s^t e^{(t-r)A}F(r, X(r, s, x)) dr + W_A(s, t), \quad t \geq 0,$$

where  $W_A(s, t)$  is the stochastic convolution

$$W_A(s, t) = \int_s^t e^{(t-r)A}(-A)^{-\gamma/2} dW(s), \quad t \geq 0.$$

(ii) Let  $x \in H$  and  $s \in [0, T]$ . We say that  $X(\cdot, s, x) \in C_W([s, T]; H)$  is a generalized solution of problem (5) if there exists a sequence  $(x_n) \subset L^{2d}(\mathcal{O})$  such that

$$\lim_{n \rightarrow \infty} x_n = x \quad \text{in } H,$$

mild solutions  $X(\cdot, s, x_n)$  exist and

$$\lim_{n \rightarrow \infty} X(\cdot, s, x_n) = X(\cdot, s, x) \quad \text{in } C_W([s, T]; H).$$

We shall denote by  $X(t, s, x)$  mild and generalized solutions of (5).

**Theorem 3.** The following statements are true.

(i) If  $x \in L^{2d}(\mathcal{O})$ , problem (5) has a unique mild solution  $X(\cdot, s, x)$ . Moreover for any  $m \in \mathbb{N}$ , there is  $c_{m,p,T} > 0$  such that

$$\mathbb{E} \left( |X(t, s, x)|_{L^{2d}(\mathcal{O})}^{2m} \right) \leq c_{m,p,T} \left( 1 + |x|_{L^{2d}(\mathcal{O})}^{2m} \right), \quad 0 \leq s \leq t \leq T.$$

(ii) If  $x \in H$ , problem (4) has a unique generalized solution  $X(\cdot, s, x)$ .

For any  $0 \leq s \leq t \leq T$ , let us consider the transition evolution operator

$$P_{s,t}\varphi(x) = \mathbb{E}[\varphi(X(t, s, x))], \quad \varphi \in C_u(H),$$

where  $X(t, s, x)$  is a generalized solution of (5). Then, given  $\nu_0 \in \mathcal{P}(H)$ , set

$$\nu_t = P_{0,t}^* \nu_0, \quad t \in [0, T]. \quad (6)$$

**Corollary 1.** Let  $m \in \mathbb{N}$  and assume that  $\nu_0 \in \mathcal{P}(H)$  satisfies

$$\int_H |x|_{L^{2d}(\mathcal{O})}^{2m} \nu_0(dx) < +\infty.$$

Then we have

$$\int_H |x|_{L^{2d}(\mathcal{O})}^{2m} \nu_t(dx) \leq c_{m,p,T} \int_H |x|_{L^{2d}(\mathcal{O})}^{2m} \nu_0(dx).$$

**Theorem 4.** Assume that  $\nu_0 \in \mathcal{P}(H)$  satisfies

$$\int_H (|x|^4 + |F(0, x)|^2 + |x|^4 |F_0(0, x)|^4) \nu_0(dx) < +\infty.$$

Then the operator  $L_0$  associated with (5) is closable and its closure is  $m$ -dissipative in  $L^2([0, +\infty) \times H, \nu)$ , where  $\nu(dt, dx) = \nu_t(dx)dt$  and  $\nu_t$  is defined by (6).

On the dissipativity of Kolmogorov operators see also [5]–[7].

This work has been supported by the RFBR projects 07-01-00536, 05-01-02941-JF, 06-01-39003 GFEN, the DFG Grant 436 RUS 113/343/0(R), the ARC Discovery Grant DP0663153 (Sydney), the programme SFB 701 at the University of Bielefeld, and the research programme “Equazioni di Kolmogorov” of the Italian “Ministero della Ricerca Scientifica e Tecnologica”.

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