# Stochastic Schrödinger equation driven by cylindrical Wiener process and fractional Brownian motion

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**Abstract.** In this paper we study the properties of the solution of a stochastic nonlinear equation of Schrödinger type, which is perturbed by a cylindrical Wiener process and an additive cylindrical fractional Brownian motion with Hurst parameter in the interval  $(\frac{1}{2}, 1)$ . The existence of the solution and the existence of the Malliavin derivative are proved.

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## 1. Introduction

In physics, specifically in quantum mechanics, the Schrödinger equation is an equation that describes how the quantum state of a physical system changes in time.

We describe the Schrödinger equation for a harmonic oscillator subject to a periodic electric field: a particle of mass m, electric charge Q, is displaced along the x-axis  $(x \in \mathbb{R})$  and subject to a force  $-m\omega_0^2 x$  (for all t > 0) and to an electric field  $E \sin(\omega t)$  directed along the x-axis

$$i\hbar\frac{\partial}{\partial t}X(x,t) = \left(-\frac{\hbar^2}{2m}\nabla^2 + \frac{1}{2}m\omega_0^2 x^2 + QEx\sin(\omega t)\right)X(x,t), \quad x \in \mathbb{R}, t > 0,$$
  
$$X(\cdot,0) = X_0$$

where *i* is the imaginary unit,  $-\frac{\hbar^2}{2m}\nabla^2$  is the kinetic energy operator,  $\hbar$  is Planck's constant, the complex valued function X is the wave function at position x at time t,  $X_0$  is the initial condition (see [8], p. 639).

Many authors investigated stochastic equations of Schrödinger type: The case of additive noise is considered in [11], [13], while the case of multiplicative noise is discussed in [2], [9], [10], [16]. In these papers the existence of a mild solution is investigated. Different approaches to linear and nonlinear stochastic Schrödinger equations perturbed by cylindrical Brownian motions are given in [14] and [15].

In this paper we study the properties of the solution of a stochastic nonlinear equation of Schrödinger type, which is perturbed by a cylindrical Wiener process and an additive cylindrical fractional Brownian motion. Consequently, this model respects as well fluctations of a Brownian motion as additive disturbances with long range dependence. This paper completes the results about stochastic equations of Schrödinger type given in [5] by considering also a cylindrical fractional Brownian motion with Hurst parameter in the interval  $(\frac{1}{2}, 1)$ . We use the framework of stochastic evolution equations driven by fractional noise developed by T.E. Duncan, B. Pasik-Duncan, B. Maslowski [12] and M. Röckner and Y. Wang [17]. The existence results are derived by using the properties of Schrödinger type equations developed in [5]. Smoothness properties such as the existence of the Malliavin derivative are also proved. The Malliavin derivatives can be used to calculate conditional expectations or chaos decompositions of stochastic processes (see [3], [7]).

This paper has the following structure: In Section 2 we introduce the list of assumptions and give the definition of the solution. In Section 3 we briefly present the two stochastic integrals that appear in the equation which is investigated. The existence of the solution is derived in Section 4. Section 5 contains results about infinite dimensional Malliavin derivatives and the existence of the Malliavin derivative of the solution is proved.

## 2. Assumptions and formulation of the problem

We consider  $(\Omega, \mathcal{F}, (\mathcal{F}_t)_{t\geq 0}, P)$  to be a filtered complete probability space. Let  $(V, (\cdot, \cdot)_V)$  and  $(H, (\cdot, \cdot))$  be separable complex Hilbert spaces, such that  $(V, H, V^*)$  forms a triplet of rigged Hilbert spaces. Let K be a separable real Hilbert space. We consider  $(W(t))_{t\geq 0}$  to be a K-valued cylindrical Wiener process adapted to the filtration  $(\mathcal{F}_t)_{t\geq 0}$  and  $(B^h(t))_{t\geq 0}$  to be a K-valued cylindrical Brownian motion with Hurst index  $h \in (\frac{1}{2}, 1)$  adapted to the filtration  $(\mathcal{F}_t)_{t\geq 0}$ .

We study the properties of the *variational solution* X of the following stochastic nonlinear evolution equation of Schrödinger type

$$(X(t), v) = (X_0, v) - i \int_0^t \langle AX(s), v \rangle ds + i \int_0^t (f(s, X(s)), v) ds \qquad (2.1)$$
$$+ i (\int_0^t g(s, X(s)) dW(s), v) + i (\int_0^t b(s) dB^h(s), v)$$

for a.e.  $\omega \in \Omega$  and all  $t \in [0, T], v \in V$ .

We assume that:

**[I]**  $X_0$  is  $\mathcal{F}_0$ -measurable,  $X_0 \in L^2(\Omega; V)$ ;

[A]  $A: V \to V^*$  has the following properties:

- A is linear and continuous  $||Au||_{V^*} \le c_A ||u||_V$  for all  $u \in V$ ;
- $\langle Au, v \rangle = \overline{\langle Av, u \rangle}$  for all  $u, v \in V$ ;
- there exists constants  $\alpha_1 \in \mathbb{R}$  and  $\alpha_2 > 0$ , such that for all  $v \in V$  it holds

$$\langle A(v), v \rangle \ge \alpha_1 \|v\|^2 + \alpha_2 \|v\|_V^2$$

• Let  $(h_n)_n \subset H$  be the eigenvectors of the operator A, for which we assume that  $Ah_n \in H$  for all  $n \in \mathbb{N}$  and  $(h_n)_n$  is a complete orthonormal system in H.

**[f]**  $f: \Omega \times [0,T] \times H \to H$  is a measurable function, which is  $\mathcal{F}_t$ -adapted for each  $t \in [0,T]$ :

(1) there exists a constant  $c_f > 0$  such that for a.e.  $\omega \in \Omega$  it holds

$$||f(t,u) - f(t,v)||^2 \le c_f ||u - v||^2$$
 for all  $t \in [0,T], u, v \in H;$ 

(2) for a.e.  $\omega \in \Omega$  and all  $t \in [0,T], u \in V$  we have  $f(t,u) \in V$  and there exists  $k_f > 0$  such that

$$||f(t,u)||_V^2 \le k_f (1+||u||_V^2);$$

**[g]**  $g : \Omega \times [0,T] \times H \to L_2(K,H)$  is a measurable function, which is  $\mathcal{F}_t$ -adapted for each  $t \in [0,T]$ :

(1) there exists a constant  $c_g > 0$  such that for a.e.  $\omega \in \Omega$  it holds

$$||g(t,u) - g(t,v)||^2_{L_2(K,H)} \le c_g ||u-v||^2$$
 for all  $t \in [0,T], u, v \in H$ ;

(2) for a.e.  $\omega \in \Omega$  and all  $t \in [0,T]$ ,  $u \in V$  we have  $g(t,u) \in L_2(K,V)$  and there exists  $k_g > 0$  such that

$$||g(t,u)||^2_{L_2(K,V)} \le k_g(1+||u||^2_V);$$

**[b]**  $b: [0,T] \to L_2(K,V)$  and for each  $u \in K$  we have  $b(\cdot)u \in L^p([0,T];V)$  for some  $p > \frac{1}{h}$  and it holds

$$\int_{0}^{T} \int_{0}^{T} \|b(r)\|_{L_{2}(K,V)} \|b(s)\|_{L_{2}(K,V)} |r-s|^{2h-2} dr ds < \infty$$

#### 3. The stochastic integrals

In this section we briefly present the definitions of the stochastic integrals we considered in (2.1). Let  $(e_n)_n$  be an orthonormal basis in K.

For the K-valued cylindrical Wiener process  $(W(t))_{t\geq 0}$  and for g:  $\Omega \times [0,T] \times H \rightarrow L_2(K,H)$  satisfying [g]-(1) the stochastic integral  $\int g(s,v)dW(s)$  ( $v \in H$  fixed) is defined as a zero mean H-valued Gauss-

ian random variable given by

$$\int_{0}^{T} g(s,v)dW(s) := \sum_{n=1}^{\infty} \int_{0}^{T} g(s,v)e_n dw_n(s),$$

where the series above converges in  $L^2(\Omega; H)$  and  $((w_n(t))_{t>0})_n$  is a sequence of mutually independent real-valued Brownian motions. One can prove that

$$E \left\| \int_{0}^{T} g(s, v) dW(s) \right\|^{2} = \sum_{n=1}^{\infty} E \left\| \int_{0}^{T} g(s, v) e_{n} dw_{n}(s) \right\|^{2}$$
$$= \sum_{n=1}^{\infty} E \int_{0}^{T} \|g(s, v) e_{n}\|^{2} ds = E \int_{0}^{T} \|g(s, v)\|^{2}_{L_{2}(K, H)} ds < \infty$$

For 0 < r < 1/(2-2h) the function  $\phi : [0,T] \to \mathbb{R}$  defined by  $\phi(u) =$  $h(2h-1)|u|^{2h-2}$  belongs to the space  $L^r([0,T];\mathbb{R})$ .

If p > 1/h, then by Theorem 3.9.4 in [4], there exists  $C_T > 0$  such that for any function  $\eta, \varphi \in L^p([0,T];\mathbb{R})$  it holds

$$\int_{0}^{T} \int_{0}^{T} |\eta(u)\varphi(v)\phi(u-v)| du dv \le C_{T} \|\varphi\|_{L^{p}([0,T];\mathbb{R})} \|\eta\|_{L^{p}([0,T];\mathbb{R})}$$

If  $(\beta^h(t))_{t\geq 0}$  is a real-valued fractional Brownian motion with Hurst index  $h \in (\frac{1}{2}, 1)$ , and  $\varphi \in L^p([0, T]; \mathbb{R})$ , then the stochastic integral  $\int_{\Omega} \varphi(s) d\beta^{h}(s) \in L^{2}(\Omega; \mathbb{R}) \text{ is defined as a zero mean real-valued Gaussian}$ 

random variable, such that

has

$$E\left(\int_{0}^{T}\varphi(s)d\beta^{h}(s)\int_{0}^{T}\varphi(s)d\beta^{h}(s)\right) = E\int_{0}^{T}\int_{0}^{T}\varphi(u)\varphi(v)\phi(u-v)dudv.$$
  
If  $\varphi \in L^{p}([0,T];\mathbb{R})$  with  $p > \frac{1}{h}$ , then the process  $\left(\int_{0}^{t}\varphi(s)d\beta^{h}(s)\right)_{t \ge 0}$   
*P*-a.s. continuous sample paths (see [18] Lemma 2.0.17).

Let  $(k_n)_n$  be an orthonormal basis in K.

For the K-valued cylindrical fractional Brownian motion  $(B^h(t))_{t>0}$  and for  $b: [0,T] \to L_2(K,V)$  satisfying assumption [b] the stochastic integral  $\int b(s)dB^{h}(s)$  is defined as a zero mean V-valued Gaussian random variable given by

$$\int_{0}^{T} b(s) dB^{h}(s) := \sum_{n=1}^{\infty} \int_{0}^{T} b(s) k_n d\beta_n^{h}(s),$$

where the series above converges in  $L^2(\Omega; V)$  and  $\left((\beta_n^h(t))_{t\geq 0}\right)_n$  is a sequence of mutually independent real-valued fractional Brownian motions each with Hurst parameter h. One can prove that

$$E \left\| \int_{0}^{T} b(s) dB^{h}(s) \right\|_{V}^{2} = \sum_{n=1}^{\infty} E \left\| \int_{0}^{T} b(s) k_{n} d\beta_{n}^{h}(s) \right\|_{V}^{2}$$
$$= \sum_{n=1}^{\infty} \int_{0}^{T} \int_{0}^{T} (b(r) k_{n}, b(s) k_{n})_{V} \phi(r, s) dr ds$$
$$\leq \int_{0}^{T} \int_{0}^{T} \|b(r)\|_{L_{2}(K,V)} \|b(s)\|_{L_{2}(K,V)} \phi(r, s) dr ds < \infty.$$

For more details see for example [12], [18].

For a.e.  $\omega \in \Omega$  and for each  $t \in [0, T]$  we denote by

$$Z(t) := \int_0^t b(s) dB^h(s),$$

which is obviously a V-valued process adapted to  $(\mathcal{F}_t)_{t\geq 0}$ .

**Proposition 3.1.** [18, Corollary 2.0.16, Lemma 2.0.17] The process  $(Z(t))_{t \in [0,T]}$  has a continuous version in V and in H and

$$E\int_0^T \|Z(s)\|_V^2 ds < \infty.$$

**Remark 3.2.** The stochastic integral Z(t) can also be represented by a stochastic integral with respect to the cylindrical Wiener process W (see [3], [6]). For  $f : \mathbb{R} \to \mathbb{C}$  and  $\frac{1}{2} < h < 1$  we introduce the operator

$$(M^h f)(x) = c_h \int_{\mathbb{R}} \frac{f(t)}{|t - x|^{3/2 - h}} dt,$$

where  $c_h = [2\Gamma(h-1/2)\cos(1/2\pi(h-1/2))]^{-1}(\Gamma(2h+1)\sin(\pi h))^{1/2}$  and f is chosen in such a manner that  $(M^h f) \in L^2(\mathbb{R})$ . If f is concentrated on [0,T], then we consider [0,T] instead of  $\mathbb{R}$ . If

$$\sum_{n=1}^{\infty}\sum_{j=1}^{\infty}\int_{0}^{T}\left(\left(M^{h}\left(b(\cdot)k_{n},h_{j}\right)\right)(s)\right)^{2}ds<\infty,$$

then

$$\int_0^t b(s)dB^h(s) = \sum_{j=1}^\infty \sum_{n=1}^\infty \int_0^t \left( M^h\left(b(\cdot)k_n, h_j\right) \right)(s)dw_n(s)h_j.$$

## 4. Existence of the solution

**Theorem 4.1.** Assume that [I], [A], [f], [g], [b] are satisfied. Equation (2.1) admits a unique solution  $X \in L^2(\Omega \times [0,T]; V) \cap L^2(\Omega; C([0,T]; H)).$ 

In order to prove the existence of the solution of (2.1), we first transform it equivalently into an equation of Schrödinger type studied in [5]. For a.e.  $\omega \in \Omega$  and for each  $t \in [0, T], v \in H$  we denote by

- U(t) := X(t) iZ(t).
- $F(\omega, t, v) := f(\omega, t, v + iZ(\omega, t)),$
- $G(\omega, t, v) := g(\omega, t, v + iZ(\omega, t)).$

Observe that for a.e.  $\omega \in \Omega$  and all  $t \in [0, T], u, v \in H$  it holds

$$||F(t,u) - F(t,v)||^2 \le c_f ||u - v||^2$$
  
|G(t,u) - G(t,v)||<sup>2</sup><sub>L2(K,H)</sub> \le c\_g ||u - v||<sup>2</sup>

and for all  $u \in V$ 

$$||F(t,u)||_{V}^{2} \leq 2k_{f}(1+||u||_{V}^{2}+||Z(t)||_{V}^{2});$$
  
$$|G(t,u)||_{L_{2}(K,V)}^{2} \leq 2k_{g}(1+||u||_{V}^{2}+||Z(t)||_{V}^{2}).$$

We rewrite (2.1) equivalently as

$$(U(t),v) = (X_0,v) - i \int_0^t \langle AU(s),v \rangle ds + i \int_0^t (F(s,U(s)),v) ds \qquad (4.1)$$
$$+i(\int_0^t G(s,U(s))dW(s),v) + i \int_0^t \langle AZ(s),v \rangle ds \text{ for all } v \in V.$$

(2.1) admits a unique solution  $X \in L^2(\Omega \times [0,T]; V) \cap L^2(\Omega; C([0,T]; H))$ if and only if (4.1) admits a unique solution  $U \in L^2(\Omega \times [0,T]; V) \cap L^2(\Omega; C([0,T]; H)).$ 

The proof of the existence of a unique solution U for (4.1) is similar to the proof of Theorem 1 in [5]. For this reason one introduces Galerkin approximations: For each  $n \in \mathbb{N}$  we consider the finite dimensional spaces  $H_n := \operatorname{sp}\{h_1, h_2, \ldots, h_n\}$  (equipped with the norm induced from H) and  $K_n := \operatorname{sp}\{e_1, e_2, \ldots, e_n\}$  (equipped with the norm induced from K). We define  $\pi_n : H \to H_n$  the orthogonal projection of H on  $H_n$  by  $\pi_n h :=$  $\sum_{j=1}^n (h, h_j)h_j$ . Let  $A_n : H_n \to H_n, F_n : \Omega \times [0, T] \times H_n \to H_n, G_n : \Omega \times [0, T] \times$  $H_n \to L(K_n, H_n)$  be defined respectively by

$$A_n u = \sum_{j=1}^n \langle Au, h_j \rangle h_j, \quad F_n(t, u) = \sum_{j=1}^n (F(t, u), h_j) h_j$$
$$G_n(t, u) v = \sum_{j=1}^n (G(t, u)v, h_j) h_i \text{ for } v \in K_n$$
$$Z_n(t) = \sum_{j=1}^n (Z(t), h_j) h_j$$

,

and we denote  $X_{0n} = \pi_n X_0$  and  $W_n(t) = \sum_{j=1}^n e_j w_j(t) \in K_n$ . For a.e.  $\omega \in \Omega$ and all  $t \in [0, T]$  and all  $j = \overline{1, n}$  we consider the finite dimensional equations corresponding to (4.1)

$$(U_{n}(t), h_{j}) = (X_{0n}, h_{j}) - i \int_{0}^{t} (A_{n}U_{n}(s), h_{j}) ds \qquad (4.2)$$
  
+ $i \int_{0}^{t} (F_{n}(s, U_{n}(s)), h_{j}) ds + i (\int_{0}^{t} G_{n}(s, U_{n}(s)) dW_{n}(s), h_{j})$   
+ $i \int_{0}^{t} (A_{n}(s)Z_{n}(s), h_{j}) ds.$ 

One can show similar as in the proof of Theorem 1 in [5] (see also Remark 3 in [5]) that for all  $t \in [0, T]$  it holds

$$\lim_{n \to \infty} E \|U_n(t) - U(t)\|^2 = 0$$

and

$$\lim_{n \to \infty} E \int_0^t \|U_n(s) - U(s)\|^2 ds = 0.$$

## 5. The existence of Malliavin derivative of the solution

We briefly present some results about infinite dimensional Malliavin derivatives: We consider the random variable Y with values in a complex Hilbert space H. Y with  $E||Y||^2 < \infty$  is called a smooth random variable and we denote  $Y \in S$ , if

$$Y = \sum_{j=1}^{n} f_j \left( \int_0^T (\gamma_{1,j}(s), dW(s))_K, \dots, \int_0^T (\gamma_{n_j,j}(s), dW(s))_K \right) h_j,$$

where  $\gamma_{1,j}, \ldots, \gamma_{n_j,j} \in L^2([0,T];K)$  for  $j = 1, \ldots, n, h_j \in H, f_j \in C^{\infty}(\mathbb{R}^{n_j})$ and  $f_j$  and all its derivatives have polynomial growth for  $j = 1, \ldots, n$ .

The Malliavin derivative  $D_t Y$ ,  $(t \in [0, T])$  of  $Y \in S$  is a random variable with values in  $L_2(K, H)$  defined by

$$D_t Y = \sum_{j=1}^n \sum_{k=1}^{n_j} \frac{\partial f_j}{\partial x_k} \left( \int_0^T \left( \gamma_{1,j}(s), dW(s) \right)_K, \dots, \int_0^T \left( \gamma_{n_j,j}(s), dW(s) \right)_K \right) \cdot h_j \otimes \gamma_{k,j}(t).$$

The Malliavin derivative  $D_t$  as defined for *H*-valued smooth random variables is closable on  $L^2(\Omega; L_2(K, H))$  (see Proposition 5.1 in [7]).

Consequently, if Y is the  $L^2(\Omega; H)$  limit of a sequence  $(Y_n)_n \subset S$  so that the sequence  $(D_tY_n)_n$  convergences in  $L^2(\Omega; L_2(K, H))$ , we can define  $D_tY$  as

$$D_t Y = \lim_{n \to \infty} D_t Y_n.$$

We use the notation H(K) for the subspace of  $L^2(\Omega; H)$ , where the derivative  $D_t$  can be defined. This subspace is a separable Hilbert space equipped with the graph norm

$$||Y||_{H(K)}^2 = E||Y||^2 + E||D_tY||_{L_2(K,H)}^2$$

The following result is known (see Lemma 5.2 in [7]):

**Lemma 5.1.** Let  $Y_n \to Y$  in  $L^2(\Omega; H)$  and suppose that there is a constant C > 0 such that for all n we have

$$E \|D_t Y\|_{L_2(K,H)}^2 < C.$$

Then, the random variable Y is in the domain H(K) of the Malliavin derivative  $D_t$ .

By using Proposition 5.2 in [7] the following chain rule holds:

**Proposition 5.2.** Let M be a further separable Hilbert space. Given a random variable  $Y \in H(K)$  and a Fréchet differentiable function  $\eta : H \to M$ . Then,

$$D_t \eta(Y) = \nabla \eta D_t Y.$$

We will use the following well-known properties of  $D_t$  (see, for example [7], [3]):

**Proposition 5.3.** (1) If Y is  $\mathcal{F}_s$ -measurable and  $Y \in H(K)$ , then  $D_t Y = 0$ a.e.  $\omega \in \Omega$  and for all t > s.

(2) Let  $a(s), s \in [0,T]$  an  $\mathcal{F}_s$ -adapted  $L_2(K,H)$ -valued process which fulfills the assumptions of the Skorochod integral definition in [7]. Then, for all r > t it holds

$$D_t \int_0^r a(s) dW(s) = a(t) + \int_t^r D_t a(s) dW(s).$$

Further in this section we assume:

- 1. The assumption in Remark 3.2 is valid for the process b.
- 2. The functions f and g are deterministic.
- 3. The functions f and g are Fréchet differentiable with respect to  $x \in H$  for all  $t \in [0, T]$  and the Fréchet derivatives  $\nabla_x f(t, x)$  and  $\nabla_x g(t, x)$  are bounded in the following sense: There exists a positive constant c such that

$$\|\nabla_x f(t,x)\|_{L(H,H)}, \|\nabla_x g(t,x)\|_{L(H,L_2(K,H))} \le c$$

for all  $t \in [0, T], x \in H$ .

4. The initial condition  $X_0$  is deterministic.

**Theorem 5.4.** There exists  $D_rU(t)$  as an  $L_2(K, H)$ -valued random variable for all  $r, t \in [0, T]$ .

*Proof.* We process the proof in two steps:

Step 1: It follows from the above assumption 3 that the functions f and g are globally Lipschitz continuous. Consequently, we can consider directly

the Galerkin equations (4.2). Similar to Remark 3 in [5] we have for the variational solution U

$$\lim_{n \to \infty} E \|U_n(t) - U(t)\|^2 = 0 \text{ and } \lim_{n \to \infty} E \int_0^t \|U_n(s) - U(s)\|^2 ds = 0 \quad (5.1)$$

for all  $t \in [0,T]$ . Equation (4.2) is an Itô equation in  $V_n$  and  $H_n$  and its solution can be approximated by the method of successive approximations

$$U_{n}^{m+1}(t) = X_{0n} - i \int_{0}^{t} A_{n} U_{n}^{m}(s) ds \qquad (5.2)$$
  
+ $i \int_{0}^{t} F_{n}(s, U_{n}^{m}(s)) ds + i \int_{0}^{t} G_{n}(s, U_{n}^{m}(s)) dW_{n}(s)$   
+ $i \int_{0}^{t} A_{n}(s) Z_{n}(s) ds.$ 

for m = 0, 1, ... with  $U^0(s) \equiv X_{0n}$ .

The finite dimensional theory shows

$$\lim_{n \to \infty} E \|U_n^m(t) - U_n(t)\|^2 = 0.$$
(5.3)

Now we calculate  $D_r U_n^{m+1}(t)$ . Since  $U_n^{m+1}$  is  $\mathcal{F}_t$ -measurable we get also the  $\mathcal{F}_r$ -measurability for  $r \geq t$ . In this case it follows from Proposition 5.3  $D_r U_n^{m+1}(t) = 0$ . We now consider r < t. Then, by Proposition 5.2, Proposition 5.3 and Remark 3.2 we get

$$D_{r}U_{n}^{m+1}(t) = -i\int_{r}^{t}A_{n}D_{r}U_{n}^{m}(s)ds \qquad (5.4)$$

$$+i\int_{r}^{t}\nabla_{x}F_{n}(s,U_{n}^{m}(s))D_{r}U_{n}^{m}(s)ds$$

$$+i\int_{r}^{t}\nabla_{x}F_{n}(s,U_{n}^{m}(s))D_{r}Z_{n}(s)ds$$

$$+i\int_{r}^{t}\nabla_{x}G_{n}(s,U_{n}^{m}(s))D_{r}U_{n}^{m}(s)dW_{n}(s)$$

$$+i\int_{r}^{t}\nabla_{x}G_{n}(s,U_{n}^{m}(s))D_{r}Z_{n}(s)dW_{n}(s)$$

$$+iG_{n}(r,U_{n}^{m}(r))+i\int_{r}^{t}A_{n}(s)D_{r}Z_{n}(s)ds$$

where  $D_r Z_n(t) : K_n \to H_n$  is the linear operator defined by

$$(D_r Z_n(t)x, y) = \left(M^h\left(b_n(\cdot)x, y\right)\right)(s).$$

 $D_r Z_n(t)$  has values in  $L(K_n, V_n)$  and  $L(K_n, H_n)$ . Since the spaces are finite dimensional, the operators are also Hilbert-Schmidt operators. If we use the energy equality in the space  $L_2(K_n, H_n)$ , then we get by the assumptions of this section and by Gronwall's lemma that there is a positive constant C with

$$E \| D_r U_n^m(t) \|_{L_2(K,H)}^2 \le C$$

for all m, r, t and fixed n, since from equation (5.3) the boundedness of  $E||U_n^m(t)||^2$  follows for all m, r, t and fixed n. The constant C does not depend on n. Then we get by Lemma 5.1, from the last inequality and from equation (5.3) that  $D_r U_n(t)$  exists and

$$E\|D_r U_n(t)\|_{L_2(K,H)}^2 \le C.$$
(5.5)

**Step 2:** Since the relations (5.5) and (5.1) hold, we can use again Lemma 5.1 and get

$$E \|D_r U(t)\|_{L_2(K,H)}^2 \le C.$$

**Theorem 5.5.** Consider that the assumptions of this section hold. Then, for t > r we have

 $D_r X(t) = D_r U(t) + i(M^h b(\cdot))(r),$ where  $(M^h b(\cdot))(r) \in L_2(K, H)$  is defined by the bilinearform  $(M^h (b(\cdot)x, y))(r)$  for all  $x \in K, x \in H.$ 

*Proof.* Theorem 5.4 shows the existence of  $D_r U(t)$  and it holds  $D_r X(t) = D_r U(t) + i D_r Z(t)$ . Since b is deterministic, we get by Proposition 5.3 and Remark 3.2 for t > r

$$D_r Z(t) = (M^h b(\cdot))(r).$$

**Remark 5.6.** The Malliavin derivative is used for example to define Skorochod integrals [12] and in the optimal control theory [1]. Optimal control problems for stochastic Schrödinger equations are under preparation.

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## References

- Anh, V.V., Grecksch, W., Yong, J., Regularity of Backward Stochastic Volterra Integral Equations, Stochastic Analysis and Applications, 29(2011), no. 1, 146– 168.
- [2] Bang, O., Christiansen, P.L., If, F., Rasmussen, K.O., Gaididei, Y.B., Temperature Effects in a Nonlinear Model of Monolayer Scheibe Aggregates, Phys. Rev. E, 49(1994), 4627-4636.
- [3] Biagini, F., Oksendal, B., Sulem, A., Wallner, N., An Introduction to White Noise Theory and Malliavin Calculus for Fractional Brownian Motion, Proc. Royal Soc. London, A, 460(2004), 347-372.
- [4] Bogachev, V.I., Measure Theory, Vol. I, Springer-Verlag, New-York, 2007.
- [5] Grecksch, W., Lisei, H., Stochastic Nonlinear Equations of Schrödinger Type, to appear in Stochastic Analysis and Applications.

- [6] Grecksch, W., Roth, C., Anh, V.V., Q-Fractional Brownian Motion in Infinite Dimensions with Applications to Fractional Black-Scholes Market, Stochastic Analysis and Applications, 27(2009), no. 1, 149-175.
- [7] Carmona, R., Tehrani, M., Interest Rate Models: an Infininite Dimensional Stochastic Analysis Perspective, Springer Verlag, Berlin - Heidelberg, 2006.
- [8] Dautray, R., Lions, J.L., Mathematical Analysis And Numerical Methods For Science And Technology, Volume 5, Evolution Problems I. Springer-Verlag, Berlin, 1992.
- [9] De Bouard, A., Debussche, A., A Stochastic Nonlinear Schrödinger Equation with Multiplicative Noise, Commun. Math. Phys., 205(1999), 161–181.
- [10] De Bouard, A., Debussche, A., A Semi-discrete Scheme for the Stochastic Nonlinear Schrödinger Equation, Numer. Math., 96(2004), 733–770.
- [11] Debussche, A., Odasso, C., Ergodicity for a Weakly Damped Stochastic Nonlinear Schrödinger Equation, J. Evol. Eq., 5(2005), 317-356.
- [12] Duncan, T.E., Maslowski, B., Pasic-Duncan, B., Fractional Brownian Motion and Stochastic Equations in Hilbert Space, Stochast. Dyn., 2(2002), 225–250.
- [13] Falkovich, G.E., Kolokolov, I., Lebedev, V., Turitsyn, S.K., Statistics of Soliton-Bearing Systems with Additive Noise, Phys. Rev. E, 63(2001).
- [14] Mora, C.M., Rebolledo, R., Regularity of Solutions to Linear Stochastic Schrödinger Equations, Infin. Dimens. Anal. Quantum Probab. Relat. Top., 10(2007), 237-259.
- [15] Mora, C.M., Rebolledo, R., Basic Properties of Nonlinear Stochastic Schrödinger Equations Driven by Brownian Motions, Ann. Appl. Probab., 18(2008), no. 2, 591-619.
- [16] Rasmussen, K.O., Gaididei, Y.B., Bang, O., Christiansen, P.L., The Influence of Noise on Critical Collapse in the Nonlinear Schrödinger Equation, Physics Letters A, 204(1995), 121–127.
- [17] Röckner, M., Wang, Y., A Note on Variational Solutions to SPDE Perturbed by Gaussian Noise in a General Class, Infin. Dimens. Anal. Quantum Probab. Relat. Top., 12(2009), no. 2, 353–358.
- [18] Wang, Y., Variational Solutions to SPDE Perturbed by a General Gaussian Noise, PhD thesis, Purdue University, 2009.

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