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Abstract

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ON THE CAUCHY PROBLEM FOR STOCHASTIC PARABOLIC EQUATIONS IN HÖLDER SPACES

KAI DU AND JIAKUN LIU

ABSTRACT. In this paper, we establish a sharp $C^{2+\alpha}$ -theory for stochastic partial differential equations of parabolic type in the whole space.

AMS Subject Classification: 35R60; 60H15

1. INTRODUCTION

In this paper, we consider the Cauchy problem for second-order stochastic partial differential equations (SPDEs) of the Itô type

(1.1)
$$du = (a^{ij}u_{x^ix^j} + b^iu_{x^i} + cu + f) dt + (\sigma^{ik}u_{x^i} + \nu^k u + g^k) dw_t^k,$$

where $\{w^k\}$ are countable independent standard Wiener processes defined on a filtered complete probability space $(\Omega, \mathcal{F}, (\mathcal{F}_t)_{t\geq 0}, \mathbb{P})$, the coefficients, free terms and the unknown function u are all random fields adapted to the filtration \mathcal{F}_t that is complete and right-continuous. Equation (1.1) has many practical applications such as in probability, engineering, and economics, and has been studied since long ago (see [Roz90]). A well-known example of (1.1) is the Zakai equation arising in the nonlinear filtering problem, see [Zak69, Roz90, Par91]. Regularity theory for equation (1.1) also play a prominent role in the study of nonlinear stochastic equations, see [Wal86, DPZ92, Kry97, Cho15] and references therein.

Denote the matrices $a = (a^{ij})$ and $\sigma = (\sigma^{ik})$. The following uniform parabolic condition is assumed throughout the paper:

(1.2)
$$\lambda I_n + \sigma \sigma^* \le 2a \le \lambda^{-1} I_n \quad \text{on } \mathbf{R}^n \times [0, \infty) \times \Omega,$$

where $\lambda > 0$ is a constant, σ^* is the transposed matrix of σ , and I_n is the $n \times n$ identity matrix.

A random field u satisfying (1.1) in the sense of Schwartz distributions is often called a weak solution of (1.1), see [Roz90]. The regularity of weak solutions in Sobolev spaces has already been investigated by many researchers. Various aspects of L^2 -theory have been obtained since 1970s, see [Par75, KR77, Roz90, DPZ92] among others. A complete L^p -theory was established by Krylov [Kry96b, Kry99] in 1990s. By Sobolev's embedding, one then has the regularity in some proper $C^{2+\alpha}$ -spaces, which however requires relatively high regularities of the given data. As an *open* problem proposed by Krylov [Kry99], one desires a sharp $C^{2+\alpha}$ -theory in the sense

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that not only that for f, g belonging to a proper space \mathcal{X} , the solution belongs to some kind of stochastic $C^{2+\alpha}$ -spaces, but also that every element of this stochastic space can be obtained as a solution for certain f, g belonging to the same \mathcal{X} .

The purpose of this paper is to establish such a sharp $C^{2+\alpha}$ -theory for equation (1.1). In order to state our result, we need to define the proper Hölder space \mathcal{X} and introduce a notion of solutions.

Definition 1. A predictable random field u is called a *quasi-classical* solution of (1.1) if

(i) for each $t \in (0, \infty)$, $u(\cdot, t)$ is a twice strongly differentiable function from \mathbf{R}^n to $L^{\gamma}_{\omega} := L^{\gamma}(\Omega, \mathcal{F}; \mathbf{R})$ for some $\gamma \geq 2$; and

(ii) for each $x \in \mathbf{R}^n$, the process $u(x, \cdot)$ is stochastically continuous and satisfies the integral equation

$$u(x, T_1) - u(x, T_0)$$

$$= \int_{T_0}^{T_1} \left[a^{ij}(x, t) u_{x^i x^j}(x, t) + b^i(x, t) u_{x^i}(x, t) + c(x, t) u(x, t) + f(x, t) \right] dt$$

$$+ \int_{T_0}^{T_1} \left[\sigma^{ik}(x, t) u_{x^i}(x, t) + \nu^k(x, t) u(x, t) + g^k(x, t) \right] dw_t^k$$

almost surely (a.s.) for all $0 \leq T_0 < T_1 < \infty$.

If furthermore, $u(\cdot, t, \omega) \in C^2(\mathbf{R}^n)$ for each $(t, \omega) \in (0, \infty) \times \Omega$, then u is a *classical* solution of (1.1).

It is well known that $L^{\gamma}_{\omega} = L^{\gamma}(\Omega, \mathcal{F}; \mathbf{R})$ is a Banach space equipped with the norm $\|\xi\|_{L^{\gamma}_{\omega}} := (\mathbb{E}|\xi|^{\gamma})^{1/\gamma}$, where $\gamma \geq 2$ is a constant fixed throughout the paper. Let T > 0 and $\mathcal{Q}_T = \mathbf{R}^n \times (0, T)$. We define the L^{γ}_{ω} -valued Hölder spaces $C^{m+\alpha}_x(\mathcal{Q}_T; L^{\gamma}_{\omega})$ and $C^{m+\alpha,\alpha/2}_{x,t}(\mathcal{Q}_T; L^{\gamma}_{\omega})$ as follows, where $\beta = (\beta_1, \ldots, \beta_n)$ denotes a multi-index and $|\beta| = \beta_1 + \cdots + \beta_n$.

Definition 2. For $m \in \mathbf{N} := \{0, 1, 2, ...\}$ and $\alpha \in (0, 1)$, the space $C_x^{m+\alpha}(\mathcal{Q}_T; L_{\omega}^{\gamma})$ consists of all predictable random fields $u : \mathcal{Q}_T \times \Omega \to \mathbf{R}$ such that $u(\cdot, t)$ is an L_{ω}^{γ} -valued strongly continuous function for each t and

(1.3)
$$|u|_{m+\alpha;\mathcal{Q}_T}^{L^{\gamma}} := |u|_{m;\mathcal{Q}_T}^{L^{\gamma}} + \max_{|\beta|=m} \sup_{t,x\neq y} \frac{\|D^{\beta}u(x,t) - D^{\beta}u(y,t)\|_{L^{\gamma}_{\omega}}}{|x-y|^{\alpha}} < \infty,$$

where $|u|_{m;Q_T}^{L_{\omega}^{\gamma}} = \max_{|\beta| \leq m} \sup_{(x,t) \in Q_T} \|D^{\beta}u(x,t)\|_{L_{\omega}^{\gamma}}$, and the derivatives are defined with respect to the spatial variable in the strong sense, see [HP57].

Using the parabolic module $|X|_{\mathbf{p}} := |x| + \sqrt{|t|}$ for $X = (x, t) \in \mathbf{R}^n \times \mathbf{R}$, we define $C_{x,t}^{m+\alpha,\alpha/2}(\mathcal{Q}_T; L_{\omega}^{\gamma})$ to be the set of all $u \in C_x^{m+\alpha}(\mathcal{Q}_T; L_{\omega}^{\gamma})$ such that

(1.4)
$$|u|_{(m+\alpha,\alpha/2);\mathcal{Q}_T}^{L^{\gamma}} := |u|_{m;\mathcal{Q}_T}^{L^{\gamma}} + \max_{|\beta|=m} \sup_{X \neq Y} \frac{\|D^{\beta}u(X) - D^{\beta}u(Y)\|_{L^{\gamma}_{\omega}}}{|X - Y|_{\mathbf{p}}^{\alpha}} < \infty.$$

Similarly, we can define the norms (1.3) and (1.4) over a domain $Q = \mathcal{O} \times I$, for any domains $\mathcal{O} \subset \mathbf{R}^n$ and $I \subset \mathbf{R}$. See §2.1 for more general definitions.

We can now state our main result, while a detailed explanation of the coefficients and free terms has to be postponed to Assumption (H) in §2.2.

Theorem 1.1. Assume that the classical C_x^{α} -norms of $a^{ij}, b^i, c, \sigma^i, \sigma_x^i, \nu, \nu_x$ are all dominated by a constant K uniformly in $(t, \omega) \in (0, T) \times \Omega$, and the condition (1.2) is satisfied. If $f \in C_x^{\alpha}(\mathcal{Q}_T; L_{\omega}^{\gamma})$, $g \in C_x^{1+\alpha}(\mathcal{Q}_T; L_{\omega}^{\gamma})$ for some $\gamma \geq 2$, then equation (1.1) with a zero initial condition admits a unique quasi-classical solution u in $C_{x,t}^{2+\alpha,\alpha/2}(\mathcal{Q}_T; L_{\omega}^{\gamma})$.

The Cauchy problem with nonzero initial value can be easily reduced to our case by a simple transform. Such an established $C^{2+\alpha}$ -theory is sharp in the sense that as proposed by Krylov in [Kry99]. We remark that by an anisotropic Kolmogorov continuity theorem (see [DKN07]), if $\gamma > n/\alpha$, the above obtained quasi-classical solution u has a $C_x^{2+\delta}$ modification with $\delta < \alpha - n/\gamma$; and if $\gamma > (n+2)/\alpha$, then uhas a $C_{x,t}^{2+\delta,\delta/2}$ modification with $\delta < \alpha - (n+2)/\gamma$.

Our result can be applied to a wide range of nonlinear filtering problems. For example, the Zakai equation is the homogeneous case of (1.1) as the terms f and g vanish, and is often associated with a deterministic initial value condition. As an application of Theorem 1.1, we have a more general result that embracing the Zakai equation.

Corollary 1.2. Under the hypotheses of Theorem 1.1, if the initial value $u(\cdot, 0) \in C^{2+\alpha}(\mathbf{R}^n)$, the free terms $f \in L^{\infty}([0,\infty); C^{\alpha}(\mathbf{R}^n))$ and $g \in L^{\infty}([0,\infty); C^{1+\alpha}(\mathbf{R}^n))$ are all nonrandom, then equation (1.1) admits a unique classical solution.

Note that in Corollary 1.2, the coefficients a, b, c, σ, ν are allowed to be random and merely required to satisfy natural regularity assumptions. To the best of our knowledge, this is a new result concerning the classical solution of the Zakai equation.

The solvability of SPDEs in L^{γ}_{ω} -valued Hölder spaces was previously studied by Rozovsky [Roz75] and Mikulevicius [Mik00]. However, they both need to assume the leading coefficient *a* is deterministic and there is no derivatives of *u* in the stochastic term, namely $\sigma \equiv 0$. Such a strong restriction excludes many interesting examples and applications. Moreover, neither of them addressed the time-continuity of secondorder derivatives of *u*, which is now obtained in our Theorem 1.1. For more related results under other appropriate assumptions, we refer the reader to, for example [Kun82, Wal86, Fun91, CJ94, BMSS95] and references therein. Most recently, Hairer [Hai14] created an abstract theory of regularity structures for SPDEs including multilevel Schauder estimates. Our approach in this paper is totally different to that of [Hai14].

The solvability in Theorem 1.1 can be derived by the standard method of continuity (see $\S2.2$), once we have the following Schauder estimate.

Theorem 1.3. Under the hypotheses of Theorem 1.1, letting u be a quasi-classical solution of (1.1) and $u(\cdot, 0) = 0$, there is a positive constant C depending only on $n, \lambda, \gamma, \alpha$ and K such that

(1.5)
$$|u|_{(2+\alpha,\alpha/2);\mathcal{Q}_T}^{L^{\omega}} \leq C e^{CT} (|f|_{\alpha;\mathcal{Q}_T}^{L^{\omega}} + |g|_{1+\alpha;\mathcal{Q}_T}^{L^{\omega}}).$$

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In non-stochastic cases, the Schauder estimate is one of the most important estimates for elliptic and parabolic equations, which was traditionally built upon the potential theory, and then was obtained via different approaches by, for instance, Campanato [Cam64], Trudinger [Tru86], Schlag [Sch96], Simon [Sim97], and others. Also, perturbation arguments were used by Safonov [Saf84], Caffarelli [Caf89], and Wang [Wan06], which can be applied to fully nonlinear equations. However, each individual method of the above has some essential defect when applied to the SPDEs, partially because of the adaptedness issues, and also the absent of a proper maximum principle for the SPDEs.¹

In our proof of Theorem 1.3, we adopt the perturbation scheme from Wang's work [Wan06], while instead of using the maximum principle, we establish certain integral-type estimates as inspired by the work of Trudinger [Tru86]. See also [DL16], where many of our results were first announced.

The paper is organised as follows. In Section 2, we introduce some notation and extend Definition 2 to general cases in §2.1, which will be used in subsequent sections. In §2.2, by assuming having Theorem 1.3, we prove Theorem 1.1 via the method of continuity. In Sections 3 and 4, we consider a model equation

(1.6)
$$du = (a^{ij}u_{x^ix^j} + f) dt + (\sigma^{ik}u_{x^i} + g^k) dw_t^k$$

where the random coefficients a and σ are independent of x. We first prove some auxiliary estimates in §3, and then establish the interior Hölder estimate in §4, which is the crucial ingredient of obtaining the Schauder estimate (1.5). In Section 5, we prove Theorem 1.3 by establishing the global Schauder estimate for the Cauchy problem of (1.1). Some properties and approximation of L^{γ}_{ω} -valued continuous functions are proved in Appendix.

2. Preliminaries

2.1. Notation. For a function u of $x = (x^1, \ldots, x^n) \in \mathbb{R}^n$, we denote

$$u_i = D_i u = u_{x^i}, \quad u_{ij} = D_{ij} u = u_{x^i x^j}, \quad Du = u_x = (u_1, \dots, u_n)$$

Hereafter, $\beta = (\beta_1, \dots, \beta_n)$ with $\beta_i \in \mathbf{N} = \{0, 1, 2, \dots\}$ is a multi-index; we denote $D^{\beta} = D^{\beta_1} \cdots D^{\beta_n}, \quad |\beta| = \beta_1 + \cdots + \beta_n.$

For $m \in \mathbf{N}$ we denote $D^m u$ the set of all *m*-order derivatives of *u*. These $D^m u(x)$ for each *x* are regarded as elements of a Euclidean space of proper dimension.

Let \mathcal{O} be a domain in \mathbb{R}^n , $I \subset \mathbb{R}$ be an interval, and $Q := \mathcal{O} \times I$. Let E be a Banach space. For a function $h : \mathcal{O} \to E$, we define

$$[h]_{0;\mathcal{O}}^{E} = |h|_{0;\mathcal{O}}^{E} := \sup_{x \in \mathcal{O}} ||h(x)||_{E}; \text{ and}$$
$$[h]_{\alpha;\mathcal{O}}^{E} := \sup_{x,y \in \mathcal{O}, \ x \neq y} \frac{||h(x) - h(y)||_{E}}{|x - y|^{\alpha}} \text{ for } \alpha \in (0, 1).$$

¹Two different types of maximum principle for SPDEs were obtained in [DMS05] and [Kry07], respectively, but neither is suitable for our circumstance.

Then for $m \in \mathbf{N}$ and $\alpha \in (0, 1)$, denote

$$|h|_{m;\mathcal{O}}^{E} := \max_{|\beta| \le m} |D^{\beta}h|_{0;\mathcal{O}}^{E},$$
$$|h|_{m+\alpha;\mathcal{O}}^{E} := |h|_{m;\mathcal{O}}^{E} + \max_{|\beta|=m} [D^{\beta}h]_{\alpha;\mathcal{O}}^{E}$$

Here and below, all the derivatives of an E-valued function are defined with respect to the spatial variable in the strong sense, see [HP57].

For a function $u: Q = \mathcal{O} \times I \to E$, we define

$$[u]_{\alpha;Q}^{E} := \sup_{t \in I} [u(\cdot, t)]_{\alpha;\mathcal{O}}^{E},$$
$$|u|_{m+\alpha;Q}^{E} := \sup_{t \in I} |u(\cdot, t)|_{m+\alpha;\mathcal{O}}^{E}$$

Letting $|X|_{\mathbf{p}} = |(x,t)|_{\mathbf{p}} = |x| + \sqrt{|t|}$ be the parabolic modulus of $X = (x,t) \in \mathbf{R}^n \times \mathbf{R}$, we further define

$$[u]_{(m+\alpha,\alpha/2);Q}^{E} := \max_{|\beta|=m} \sup_{X,Y \in Q, \ X \neq Y} \frac{\|D_{x}^{\beta}u(X) - D_{x}^{\beta}u(Y)\|_{E}}{|X - Y|_{\mathbf{p}}^{\alpha}},$$

$$|u|_{(m+\alpha,\alpha/2);Q}^{E} := |u|_{m;Q}^{E} + [u]_{(m+\alpha,\alpha/2);Q}^{E}.$$

In the following context, the space E is either a Euclidean space or L^{γ}_{ω} , where $\gamma \in [2, \infty)$ is a fixed constant. We omit the superscript when E is a Euclidean space. In the case of $E = L^{\gamma}_{\omega}$, we introduce some new notation:

$$] \cdot \llbracket_{\dots} := | \cdot |_{\dots}^{L^{\gamma}_{\omega}}, \quad \llbracket \cdot \rrbracket_{\dots} := [\cdot]_{\dots}^{L^{\gamma}_{\omega}}.$$

For instance, $[\![u]\!]_{m+\alpha;Q} = |u|^{L^{\gamma}_{\omega}}_{m+\alpha;Q}$, and $[\![u]\!]_{(\alpha,\alpha/2);Q} = [u]^{L^{\gamma}_{\omega}}_{(\alpha,\alpha/2);Q}$. Using the above notation, the spaces $C^{m+\alpha}_x(Q; L^{\gamma}_{\omega})$ and $C^{m+\alpha,\alpha/2}_{x,t}(Q; L^{\gamma}_{\omega})$ defined

Using the above notation, the spaces $C_x^{m+\alpha}(Q; L_{\omega}^{\gamma})$ and $C_{x,t}^{m+\alpha,\alpha/2}(Q; L_{\omega}^{\gamma})$ defined in Definition 2 are the sets of all predictable random fields $u: Q \times \Omega \to \mathbf{R}$ such that $[\![u]\!]_{m+\alpha;Q}$ and $[\![u]\!]_{(m+\alpha,\alpha/2);Q}$ are finite, respectively.

2.2. The solvability. Let L and A^k be differential operators

 $L = a^{ij}D_{ij} + b^iD_i + c, \quad \text{and} \quad \Lambda^k = \sigma^{ik}D_i + \nu^k.$

The Cauchy problem under consideration can be written as

(2.1)
$$\begin{aligned} \mathrm{d}u &= (Lu+f)\,\mathrm{d}t + (\Lambda^k u + g^k)\,\mathrm{d}w_t^k \quad \text{in } \mathcal{Q} := \mathbf{R}^n \times [0,\infty), \\ u(\cdot,0) &= 0 \quad \text{in } \mathbf{R}^n. \end{aligned}$$

Throughout the paper, we assume that

(H) For all i, j = 1, ..., n, the random fields a^{ij}, b^i, c and f are real-valued, and σ^i, ν and g are ℓ^2 -valued; all of them are predictable. a^{ij} and σ^i satisfy the stochastic parabolic condition (1.2). For some $\alpha \in (0, 1)$ there exists a constant K such that $\max\{|a^{ij}|_{\alpha;\mathcal{Q}}, |b^i|_{\alpha;\mathcal{Q}}, |c|_{\alpha;\mathcal{Q}}, |\sigma^i|_{1+\alpha;\mathcal{Q}}, |\nu|_{1+\alpha;\mathcal{Q}}\} \leq K$ for all $\omega \in \Omega$.

Recall that $Q_T = \mathbf{R}^n \times (0, T)$, and T > 0. Using the notation in §2.1, the Schauder estimate in Theorem 1.3 can be written as: There is a positive constant C depending only on $n, \lambda, \gamma, \alpha$ and K, such that

(2.2)
$$\|u\|_{(2+\alpha,\alpha/2);\mathcal{Q}_T} \le C e^{CT} (\|f\|_{\alpha;\mathcal{Q}_T} + \|g\|_{1+\alpha;\mathcal{Q}_T})$$

for any T > 0, where u is a quasi-classical solution of the Cauchy problem (2.1).

Proof of Theorem 1.1. With the above a priori estimates in hand, we can obtain the solvability of the Cauchy problem (2.1) by the method of continuity. Consider

(2.3)
$$du = (L_s u + f) dt + (\Lambda_s^k u + g^k) dw_t^k, \quad u(\cdot, 0) = 0,$$

where $s \in [0, 1]$ and

$$L_s := sL + (1-s)\Delta, \quad \Lambda_s^k := s\Lambda^k$$

Evidently, the solutions of (2.3) satisfy the a priori estimate (2.2) with the constant C independent of s. In view of [GT01, Theorem 5.2], it suffices to show the solvability of the stochastic heat equation (the case s = 0):

(2.4)
$$\operatorname{d} u = (\Delta u + f) \operatorname{d} t + g^k \operatorname{d} w_t^k, \quad u(\cdot, 0) = 0.$$

Set $f^{\varepsilon} = \varphi^{\varepsilon} * f$ and $g^{\varepsilon} = \varphi^{\varepsilon} * g$, where $\varphi^{\varepsilon}(x) = \varepsilon^n \varphi(x/\varepsilon)$ and φ is a nonnegative and symmetric mollifier defined on \mathbf{R}^n (see Appendix). Then (from Lemma A.6 in Appendix) we have that $f^{\varepsilon} \in C_x^{\alpha}(\mathcal{Q}; L_{\omega}^{\gamma})$ and $g^{\varepsilon} \in C_x^{1+\alpha}(\mathcal{Q}; L_{\omega}^{\gamma})$ satisfying

(2.5)
$$[f^{\varepsilon} - f]_{\alpha/2;\mathcal{Q}_T} + [g^{\varepsilon} - g]_{1+\alpha/2;\mathcal{Q}_T} \to 0, \quad \text{as } \varepsilon \to 0.$$

Moreover, $f^{\varepsilon}(x,t,\omega)$ and $g^{\varepsilon}(x,t,\omega)$ are smooth in x for any (t,ω) , and $f^{\varepsilon}, g^{\varepsilon} \in C^{m}(\mathcal{Q}_{T}; L^{\gamma}_{\omega})$ for all $m \in \mathbf{N}$, so by Fubini's theorem,

$$\mathbb{E} \int_{\mathcal{Q}_T} (1+|x|^2)^{-p} (|D^m f^{\varepsilon}(x,t)|^{\gamma} + |D^m g^{\varepsilon}(x,t)|^{\gamma}) \, \mathrm{d}x \mathrm{d}t$$

$$\leq \int_{\mathcal{Q}_T} (1+|x|^2)^{-p} (\mathbb{E} |D^m f^{\varepsilon}(x,t)|^{\gamma} + \mathbb{E} |D^m g^{\varepsilon}(x,t)|^{\gamma}) \, \mathrm{d}x \mathrm{d}t$$

$$\leq C(n,p)T ([]f^{\varepsilon} []^{\gamma}_{m;\mathcal{Q}_T} + []g^{\varepsilon} []^{\gamma}_{m;\mathcal{Q}_T}) < \infty \quad \forall m \in \mathbf{N}$$

with 2p > n. Therefore, it follows from Krylov–Rozovsky [KR82, Theorem 2.2] that (2.4) with free terms f^{ε} and g^{ε} admits a unique weak solution u^{ε} satisfying

$$\mathbb{E}\sup_{t\in[0,T]}\int_{\mathbf{R}^n}(1+|x|^2)^{-p}|D^m u^{\varepsilon}(x,t)|^{\gamma}\,\mathrm{d} x<\infty\quad\forall\,m\in\mathbf{N}$$

and by Sobolev's embedding, u^{ε} is smooth in x, and $\mathbb{E}|D^{m}u^{\varepsilon}(x,t)|^{\gamma} < \infty$ for each $(x,t) \in \mathcal{Q}_{T}$ and $m \in \mathbb{N}$ (see Lemma A.1 in Appendix). From estimate (2.2) (with $\alpha/2$ instead of α) and keeping (2.5) in mind, we have

$$[\!]u^{\varepsilon} - u^{\varepsilon'}]\!]_{2;\mathcal{Q}_T} \le C([\!]f^{\varepsilon} - f^{\varepsilon'}]\!]_{\alpha/2;\mathcal{Q}_T} + [\!]g^{\varepsilon} - g^{\varepsilon'}]\!]_{1+\alpha/2;\mathcal{Q}_T}) \to 0$$

as $\varepsilon, \varepsilon' \to 0$. Hence, u^{ε} converges to a function $u \in C^{2,0}_{x,t}(\mathcal{Q}_T; L^{\gamma}_{\omega})$ that apparently solves (2.4). The regularity and the uniqueness follow directly from the estimate (2.2).

3. AUXILIARY ESTIMATES FOR THE MODEL EQUATION

Given countable independent Wiener processes (still denoted as $\{w^k\}$) starting at time t = -1, and predictable processes $a^{ij} = \{a^{ij}(t) : t \ge -1\}$ and $\sigma^{ik} = \{\sigma^{ik}(t) : t \ge -1\}$ satisfying the stochastic parabolic condition (1.2), we consider the *model* equation

(3.1)
$$du(x,t) = [a^{ij}(t)u_{ij}(x,t) + f(x,t)] dt + [\sigma^{ik}(t)u_i(x,t) + g^k(x,t)] dw_t^k$$

with $(x,t) \in \mathbf{R}^n \times [-1,\infty)$.

We first prove some auxiliary estimates for the model equation in this section, and then proceed to the interior Hölder estimate in the next section.

Let $\mathcal{O} \subset \mathbf{R}^n$, and $H^m(\mathcal{O}) = W^{m,2}(\mathcal{O})$ be the usual Sobolev spaces. Let $I \subset \mathbf{R}$ and $Q = \mathcal{O} \times I$. Define

$$L^p_{\omega}L^q_tH^m_x(Q) := L^p(\Omega; L^q(I; H^m(\mathcal{O}))), \quad \text{for} \ p, q \in [1, \infty].$$

For r > 0, we denote

(3.2)
$$B_r(x) = \{ y \in \mathbf{R}^n : |y - x| < r \}, \quad Q_r(x, t) = B_r(x) \times (t - r^2, t],$$

and simply write $B_r = B_r(0), Q_r = Q_r(0, 0).$

Proposition 3.1. Let *m* be a positive integer, $r \in (0,1]$ and $\theta \in (0,1)$. Let $u \in L^p_{\omega}L^2_t H^{m+1}_x(Q_r)$ solve (3.1) in Q_r with $f \in L^p_{\omega}L^2_t H^{m-1}_x(Q_r)$ and $g \in L^p_{\omega}L^2_t H^m_x(Q_r)$. Then there exists a constant *C* depending only on n, p, λ, m and θ such that

$$\begin{split} \|D^{m}u\|_{L^{p}_{\omega}L^{\infty}_{t}L^{2}_{x}(Q_{\theta r})} + \|D^{m}u_{x}\|_{L^{p}_{\omega}L^{2}_{t}L^{2}_{x}(Q_{\theta r})} &\leq Cr^{-m-1}\|u\|_{L^{p}_{\omega}L^{2}_{t}L^{2}_{x}(Q_{r})} \\ + C\sum_{k=0}^{m-1}r^{-m+k+1}\|D^{k}f\|_{L^{p}_{\omega}L^{2}_{t}L^{2}_{x}(Q_{r})} + C\sum_{k=0}^{m}r^{-m+k}\|D^{k}g\|_{L^{p}_{\omega}L^{2}_{t}L^{2}_{x}(Q_{r})}. \end{split}$$

Consequently, for $2(m - |\beta|) > n$,

$$r^{\frac{n}{2}-m+|\beta|} \|\sup_{Q_{\theta r}} |D^{\beta}u| \|_{L^{p}_{\omega}} \leq Cr^{-m-1} \|u\|_{L^{p}_{\omega}L^{2}_{t}L^{2}_{x}(Q_{r})}$$
$$+ C\sum_{k=0}^{m-1} r^{-m+k+1} \|D^{k}f\|_{L^{p}_{\omega}L^{2}_{t}L^{2}_{x}(Q_{r})} + C\sum_{k=0}^{m} r^{-m+k} \|D^{k}g\|_{L^{p}_{\omega}L^{2}_{t}L^{2}_{x}(Q_{r})}$$

where the constant C further depends on $|\beta|$.

In order to establish the above local estimates, we first show the following mixednorm estimates for the model equation (3.1).

Lemma 3.2. Let $\mathcal{Q}_T = \mathbf{R}^n \times [0, T]$, $p \geq 2$ and $m \in \mathbf{N}$. Suppose $f \in L^p_{\omega} L^2_t H^{m-1}_x$ and $g \in L^p_{\omega} L^2_t H^m_x$. Then equation (3.1) with condition u(x, 0) = 0 for all $x \in \mathbf{R}^n$ admits a unique weak solution $u \in L^p_{\omega} L^\infty_t H^m_x(\mathcal{Q}_T) \cap L^p_{\omega} L^2_t H^{m+1}_x(\mathcal{Q}_T)$, and for any multi-index β such that $|\beta| \leq m$,

$$(3.3) \|D^{\beta}u\|_{L^{p}_{\omega}L^{\infty}_{t}L^{2}_{x}} + \|D^{\beta}u_{x}\|_{L^{p}_{\omega}L^{2}_{t}L^{2}_{x}} \le C(\|D^{\beta}f\|_{L^{p}_{\omega}L^{2}_{t}H^{-1}_{x}} + \|D^{\beta}g\|_{L^{p}_{\omega}L^{2}_{t}L^{2}_{x}}).$$

where $C = C(n, p, T, \lambda)$.

Proof. The special case of p = 2 follows from the L^2 -theory of SPDEs, for instance, see [Roz90, Theorem 4.1.2]. We prove the general cases of $p \ge 2$ by induction of m. Thus it suffices to show the estimate (3.3) for m = 0, namely

(3.4)
$$\|u\|_{L^{p}_{\omega}L^{\infty}_{t}L^{2}_{x}} + \|u\|_{L^{p}_{\omega}L^{2}_{t}H^{1}_{x}} \leq C(\|f\|_{L^{p}_{\omega}L^{2}_{t}H^{-1}_{x}} + \|g\|_{L^{p}_{\omega}L^{2}_{t}L^{2}_{x}}).$$

Take a stopping time $\tau: \Omega \to [0,T]$ such that

(3.5)
$$\mathbb{E}\left[\left(\sup_{t\in[0,\tau]}\int_{\mathbf{R}^n}|u(x,t)|^2\,\mathrm{d}x+\int_0^\tau\!\!\int_{\mathbf{R}^n}|u_x(x,t)|^2\,\mathrm{d}x\mathrm{d}t\right)^{\frac{p}{2}}\right]<\infty.$$

By the Itô formula (cf. [Roz90, Theorem 4.2.2]) and integration by parts,

$$\|u(t)\|_{L_x^2}^2 = \int_0^t \int_{\mathbf{R}^n} \left[-2a^{ij}u_i u_j + \|\sigma^i u_i + g\|_{\ell^2}^2 \right] dxds + 2\int_0^t \langle u(t), f(t) \rangle dt + 2\int_0^t \int_{\mathbf{R}^n} ug^k dxdw_s^k$$

where $\langle \cdot, \cdot \rangle$ denotes the duality product between and H_x^1 and H_x^{-1} . By the parabolic condition (1.2) and using Young's inequality, we have

$$\begin{split} \sup_{t\in[0,\tau]} \|u(t)\|_{L^2_x}^2 &+ \int_0^\tau \|u(t)\|_{H^1_x}^2 \,\mathrm{d}t \\ &\leq C \int_0^\tau \big[\|f(t)\|_{H^{-1}_x}^2 + \|g(t)\|_{L^2_x}^2 \big] \mathrm{d}t + C \left| \sup_{t\in[0,\tau]} \int_0^t \!\!\!\int_{\mathbf{R}^n} ug^k \,\mathrm{d}x \mathrm{d}w_s^k \right|. \end{split}$$

Then computing $\mathbb{E}[\,\cdot\,]^{p/2}$ on both sides gives us that

(3.6)
$$\mathbb{E}\left[\left(\sup_{t\in[0,\tau]}\|u(t)\|_{L^{2}_{x}}^{2}+\int_{0}^{\tau}\|u(t)\|_{H^{1}_{x}}^{2}\,\mathrm{d}t\right)^{\frac{p}{2}}\right] \\ \leq C\left(\|f\|_{L^{p}_{\omega}L^{2}_{t}H^{-1}_{x}(\mathcal{Q}_{T})}^{p}+\|g\|_{L^{p}_{\omega}L^{2}_{t}L^{2}_{x}(\mathcal{Q}_{T})}^{p}\right)+C\mathbb{E}\left[\sup_{t\in[0,\tau]}\left|\int_{0}^{t}\int_{\mathbf{R}^{n}}ug^{k}\,\mathrm{d}x\mathrm{d}w_{s}^{k}\right|^{\frac{p}{2}}\right],$$

where $C = C(\lambda, T)$. By the Burkholder–Davis–Gundy inequality (see [RY99]), the last term is dominated by

$$C\mathbb{E}\left[\left(\int_{0}^{\tau}\sum_{k=1}^{\infty}\left|\int_{\mathbf{R}^{n}}ug^{k}\,\mathrm{d}x\right|^{2}\mathrm{d}t\right)^{\frac{p}{4}}\right]$$

$$(3.7) \qquad \leq C\mathbb{E}\left[\left(\sup_{t\in[0,\tau]}\int_{\mathbf{R}^{n}}|u(x,t)|^{2}\,\mathrm{d}x\right)^{\frac{p}{4}}\left(\int_{0}^{\tau}\int_{\mathbf{R}^{n}}\|g(x,t)\|_{\ell^{2}}^{2}\,\mathrm{d}x\mathrm{d}t\right)^{\frac{p}{4}}\right]$$

$$\leq \varepsilon\mathbb{E}\sup_{t\in[0,\tau]}\|u(t)\|_{L^{2}_{x}}^{p}+C\varepsilon^{-1}\|g\|_{L^{p}_{\omega}L^{2}_{t}L^{2}_{x}(\mathcal{Q}_{T})}^{p}.$$

Taking the positive number ε sufficiently small and combining (3.7) along with (3.6), we thus obtain that

$$\mathbb{E}\bigg[\bigg(\sup_{t\in[0,\tau]}\|u(t)\|_{L^2_x}^2 + \int_0^\tau \|u(t)\|_{H^1_x}^2 \,\mathrm{d}t\bigg)^{\frac{p}{2}}\bigg] \le C\big(\|f\|_{L^p_\omega L^2_t H^{-1}_x(\mathcal{Q}_T)}^p + \|g\|_{L^p_\omega L^2_t L^2_x(\mathcal{Q}_T)}^p\big),$$

where the constant C depends only on λ and T. Then (3.4) follows by applying the above estimate to the following sequence of stopping times

$$\tau_k := \inf\left\{t \ge 0 : \sup_{s \in [0,t]} \|u(s)\|_{L^2_x}^2 + \int_0^t \|u(s)\|_{H^1_x}^2 \,\mathrm{d}s > k\right\} \wedge T,$$

and sending k to infinity. For $m \ge 1$, one can easily apply the induction argument to conclude the lemma.

Proof of Proposition 3.1. Now we are ready to prove the estimates in Proposition 3.1. It suffices to consider the case r = 1. For general r > 0, we can apply the obtained estimates for r = 1 to the rescaled function

$$v(x,t) := u(rx, r^2t), \quad \forall (x,t) \in \mathbf{R}^n \times [-1,\infty)$$

that solves the equation

(3.8)
$$dv(x,t) = [a^{ij}(r^2t)v_{ij}(x,t) + F(x,t)]dt + [\sigma^{ik}(r^2t)v_i(x,t) + G^k(x,t)]d\beta_t^k,$$

with

$$F(x,t) = r^2 f(rx, r^2 t), \quad G(x,t) = rg(rx, r^2 t), \quad \beta_t^k = r^{-1} w_{r^2 t}^k,$$

and obviously, β^k are mutually independent Wiener processes.

By induction, we shall only consider the case of m = 1. Let $\zeta \in C_0^{\infty}(\mathbb{R}^{n+1})$ be a nonnegative function such that $\zeta(x,t) = 1$ if $|(x,t)|_{\mathbb{P}} \leq \sqrt{\theta}$, where $\theta \in (0,1)$, and $\zeta(x,t) = 0$ if $|(x,t)|_{\mathbb{P}} \geq (1+\sqrt{\theta})/2$. Then $v = \zeta u$ satisfies

(3.9)
$$\mathrm{d}v = (a^{ij}D_{ij}v + \tilde{f})\,\mathrm{d}t + (\sigma^{ik}D_iv + \tilde{g}^k)\,\mathrm{d}w_t^k,$$

where

$$\tilde{f} = \zeta f - 2a^{ij}(\zeta_i u)_j + a^{ij}\zeta_{ij}u + \zeta_t u, \quad \tilde{g}^k = \zeta g^k - \sigma^{ik}\zeta_i u.$$

Applying Lemma 3.2 to (3.9) with $|\beta| = 0$, we have

$$(3.10) \|u\|_{L^p_{\omega}L^\infty_t L^2_x(Q_{\sqrt{\theta}})} + \|Du\|_{L^p_{\omega}L^2_t L^2_x(Q_{\sqrt{\theta}})} \\ \leq C \big(\|u\|_{L^p_{\omega}L^2_t L^2_x(Q_1)} + \|f\|_{L^p_{\omega}L^2_t L^2_x(Q_1)} + \|g\|_{L^p_{\omega}L^2_t L^2_x(Q_1)} \big).$$

While by choosing another cut-off function ζ such that $\zeta(x,t) = 1$ if $|(x,t)|_{\mathbf{p}} \leq \theta$, and $\zeta(x,t) = 0$ if $|(x,t)|_{\mathbf{p}} \geq \sqrt{\theta}$, and again applying Lemma 3.2 with $|\beta| = 1$, we have

$$(3.11) \|Du\|_{L^p_{\omega}L^\infty_t L^2_x(Q_{\theta})} + \|D^2u\|_{L^p_{\omega}L^2_t L^2_x(Q_{\theta})} \\ \leq C \big(\|Du\|_{L^p_{\omega}L^2_t L^2_x(Q_{\sqrt{\theta}})} + \|f\|_{L^p_{\omega}L^2_t L^2_x(Q_1)} + \|g\|_{L^p_{\omega}L^2_t H^1_x(Q_1)} \big).$$

Combining (3.10) and (3.11), the first inequality in Proposition 3.1 is proved.

In view of Sobolev's embedding theorem, the second inequality in Proposition 3.1 follows directly from the first one. In fact, from Sobolev's theorem, one has $H_x^m(Q_\theta) \subset C_x^j(Q_\theta)$ if 2(m-j) > n, see [AF03, Theorem 4.12]. Hence, $L_{\omega}^p L_t^{\infty} H_x^m(Q_\theta) \subset L_{\omega}^p L_t^{\infty} C_x^j(Q_\theta)$. More specifically, if $2(m - |\beta|) > n$, then there is a constant C depending only on $n, m, |\beta|$ and θ such that

$$\|\sup_{Q_{\theta}}|D^{\beta}u|\|_{L^{p}_{\omega}} \leq C \|u\|_{L^{p}_{\omega}L^{\infty}_{t}H^{m}_{x}(Q_{\theta})},$$

which along with (3.10) and the first inequality in Proposition 3.1 yields the desired estimate. The proof is complete.

As an immediate application, we give an estimate for equation (3.1) with the Dirichlet boundary conditions

(3.12)
$$u(0, \cdot) = 0, \quad u|_{\partial B_r} = 0.$$

Proposition 3.3. Let f and g be in $L^{\gamma}_{\omega}L^2_t H^m_x(B_r \times (0, r^2))$ for all $m \in \mathbb{N}$. Then the Dirichlet problem (3.1) and (3.12) has a unique weak solution $u \in L^2_{\omega}L^2_t H^1_x(B_r \times (0, r^2))$, and for each $t \in (0, r^2)$, $u(t, \cdot) \in L^{\gamma}(\Omega; C^m(B_{\varepsilon}))$ for all $m \ge 0$ and $\varepsilon \in (0, r)$. Moreover, there is a constant $C = C(n, \gamma)$ such that

$$(3.13) \|u\|_{L^{\gamma}_{\omega}L^{2}_{t}L^{2}_{x}(B_{r}\times(0,r^{2}))} \leq C(r^{2}\|f\|_{L^{\gamma}_{\omega}L^{2}_{t}L^{2}_{x}(B_{r}\times(0,r^{2}))} + r\|g\|_{L^{\gamma}_{\omega}L^{2}_{t}L^{2}_{x}(B_{r}\times(0,r^{2}))}).$$

Proof. The existence, uniqueness and smoothness of the weak solution of the Dirichlet problem (3.1) and (3.12) follow from [Kry94, Theorem 2.1]. And the estimate (3.13) can be derived analogously to that of (3.4) by means of Itô's formula and rescaling.

4. Interior Hölder estimates for the model equation

We continue the investigation to the model equation (3.1) with $(x,t) \in \mathbf{R}^n \times [-1,\infty)$. The aim of this section is to prove the interior Hölder estimates for (3.1). To be more general, we assume that $f \in C_x^0(\mathbf{R}^n \times [-1,\infty); L_{\omega}^{\gamma})$ and $g \in C_x^1(\mathbf{R}^n \times [-1,\infty); L_{\omega}^{\gamma})$, and f(x,t) and $g_x(x,t)$ are Dini continuous with respect to x uniformly in t, namely, the modulus of continuity defined by

$$\varpi(r) = \sup_{t \ge -1, \, |x-y| \le r} (\|f(x,t) - f(y,t)\|_{L^{\gamma}_{\omega}} + \|g_x(x,t) - g_x(y,t)\|_{L^{\gamma}_{\omega}})$$

satisfies that

$$\int_0^1 \frac{\varpi(r)}{r} \,\mathrm{d}r < \infty.$$

Recall the notation B_r , Q_r and $Q_r(x,t)$ defined in (3.2). The main estimate is the following

Theorem 4.1. Let $Z = (z, s) \in \mathbf{R}^n \times [0, \infty)$ and u be a quasi-classical solution to (3.1) in $Q_1(Z)$. Under the above settings, there is a positive constant C, depending only on n, λ and γ , such that for any $X, Y \in Q_{1/4}(Z)$,

(4.1)
$$\|u_{xx}(X) - u_{xx}(Y)\|_{L^{\gamma}_{\omega}} \leq C \left[\delta M_1 + \int_0^{\delta} \frac{\overline{\omega}(r)}{r} \,\mathrm{d}r + \delta \int_{\delta}^1 \frac{\overline{\omega}(r)}{r^2} \,\mathrm{d}r\right],$$

where $\delta = |X - Y|_{\mathbf{p}}$ and $M_1 = [[u]]_{0;Q_1(Z)} + [[f]]_{0;Q_1(Z)} + [[g]]_{1;Q_1(Z)}$.

An immediate consequence is the following interior Hölder estimate for (3.1), where we denote $Q_{r,T} = B_r \times [0,T]$ for r, T > 0. **Corollary 4.2.** Let u be a quasi-classical solution of (3.1), and $\alpha \in (0,1)$. If u(x,0) = 0 for all $x \in \mathbf{R}^n$, then there is a positive constant C, depending only on n, λ, γ and α , such that

(4.2)
$$[\![u_{xx}]\!]_{(\alpha,\alpha/2);\mathcal{Q}_{1/4,T}} \leq C \left[[\![u]\!]_{0;\mathcal{Q}_{1,T}} + \frac{[\![f]\!]_{\alpha;\mathcal{Q}_{1,T}} + [\![g]\!]_{1+\alpha;\mathcal{Q}_{1,T}}}{\alpha(1-\alpha)} \right]$$

for any T > 0, provided the right-hand side is finite.

Proof. We define $\tilde{u}(x,t), \tilde{f}(x,t)$ and $\tilde{g}(x,t)$ to be zero whenever $t \in [-1,0)$, and be equal to u(x,t), f(x,t) and g(x,t), respectively, whenever $t \ge 0$. Then it is easily verified that \tilde{u} is a quasi-classical solution (see Definition 2) of the following equation

$$d\tilde{u} = (a^{ij}\tilde{u}_{ij} + \tilde{f}) dt + (\sigma^{ik}\tilde{u}_i + \tilde{g}^k) dw_t^k$$

in the area $\mathbb{R}^n \times [-1, \infty)$. Applying (4.1) to the above equation we have

(4.3)
$$[\![\tilde{u}_{xx}]\!]_{(\alpha,\alpha/2);Q_{1/4}(X)} \leq C \left[[\![\tilde{u}]\!]_{0;Q_1(X)} + \frac{[\![\tilde{f}]\!]_{\alpha;Q_1(X)} + [\![\tilde{g}]\!]_{1+\alpha;Q_1(X)}}{\alpha(1-\alpha)} \right]$$

for any $X = (x, t) \in \mathbf{R}^n \times [0, \infty)$. Then we fix x = 0 and let t run through [0, T]; keeping in mind that \tilde{u} , \tilde{f} and \tilde{g} vanish when $t \leq 0$, and using the localization property of Hölder norms (see [Kry96a, Lemma 4.1.1]), we have

$$\begin{split} \|u_{xx}\|_{(\alpha,\alpha/2);\mathcal{Q}_{1/4,T}} &= \|\tilde{u}_{xx}\|_{(\alpha,\alpha/2);\mathcal{Q}_{1/4,T}} \\ &\leq C(n,\alpha) \sup_{0 \leq t \leq T} \left(\|\tilde{u}_{xx}\|_{(\alpha,\alpha/2);\mathcal{Q}_{1/4}(0,t)} + \|\tilde{u}\|_{0;\mathcal{Q}_{1/4}(0,t)} \right) \\ &\leq C \sup_{0 \leq t \leq T} \left[\|\tilde{u}\|_{0;\mathcal{Q}_{1}(0,t)} + \frac{\|\tilde{f}\|_{\alpha;\mathcal{Q}_{1}(0,t)} + \|\tilde{g}\|_{1+\alpha;\mathcal{Q}_{1}(0,t)}}{\alpha(1-\alpha)} \right] \\ &\leq C \left[\|u\|_{0;\mathcal{Q}_{1,T}} + \frac{\|f\|_{\alpha;\mathcal{Q}_{1,T}} + \|g\|_{1+\alpha;\mathcal{Q}_{1,T}}}{\alpha(1-\alpha)} \right]. \end{split}$$

The proof is complete.

Proof of Theorem 4.1. By means of translation, we may suppose Z = (0,0) without loss of generality.

Letting $\varphi : \mathbf{R}^n \to \mathbf{R}$ be a nonnegative and symmetric mollifier (see Appendix) and $\varphi^{\varepsilon}(x) = \varepsilon^n \varphi(x/\varepsilon)$, we define $u^{\varepsilon} = \varphi^{\varepsilon} * u$, $f^{\varepsilon} = \varphi^{\varepsilon} * f$ and $g^{\varepsilon} = \varphi^{\varepsilon} * g$. Under the condition of Theorem 4.1, it follows from Corollary A.5 (see Appendix) that

$$\begin{split} \|f^{\varepsilon} - f\|_{0;\mathbf{R}^{n}} + \|g^{\varepsilon} - g\|_{1;\mathbf{R}^{n}} \to 0, \\ \|D^{2}u^{\varepsilon}(X) - D^{2}u(X)\|_{L^{\gamma}_{\omega}} \to 0 \quad \forall X \in \mathbf{R}^{n} \times \mathbf{R}, \end{split}$$

as $\varepsilon \to 0$. Evidently, f^{ε} and Dg^{ε} are also Dini continuous and has the same modulus of continuity ϖ with f and Dg. On the other hand, from Fubini's theorem one can check that u^{ε} satisfies the model equation (3.1) in the classical sense with free terms f^{ε} and g^{ε} . Therefore, it suffices to prove the theorem for the mollified functions, and the general case is straightforward by passing the limits. Based on the above analysis and the property of mollified functions (see Lemmas A.1 and A.2 and Remark A.1 in Appendix), we may assume that f and g satisfy the following additional condition:

(A)
$$f, g \in L^{\gamma}_{\omega}L^2_t H^k_x(Q_R) \cap C^k_x(Q_R; L^{\gamma}_{\omega})$$
 for all $k \in \mathbb{N}$ and $R > 0$.

From the definition of ϖ , one can see that for any $x, y \in \mathbf{R}^n$ and $t \in \mathbf{R}$,

(4.4)
$$\|f(x,t) - f(y,t)\|_{L^{\gamma}_{\omega}} + \|g_x(x,t) - g_x(y,t)\|_{L^{\gamma}_{\omega}} \le \varpi(|x-y|), \\ \|g(x,t) - g(y,t) - g_x(x,t) \cdot (x-y)\|_{L^{\gamma}_{\omega}} \le |x-y|\varpi(|x-y|).$$

With $\rho = 1/2$, we denote

$$Q^{\kappa} = Q_{\rho^{\kappa}} = Q_{\rho^{\kappa}}(0,0), \quad \kappa = 0, 1, 2, \cdots.$$

Let us introduce the following boundary problems:

$$\begin{split} \mathrm{d} u^{\kappa} &= \left[a^{ij} u^{\kappa}_{ij} + f(0,t)\right] \mathrm{d} t + \left[\sigma^{ik} u^{\kappa}_{i} + g^{k}(0,t) + g^{k}_{x}(0,t) \cdot x\right] \mathrm{d} w^{k}_{t} \quad \text{in } Q^{\kappa}, \\ u^{\kappa} &= u \quad \text{on } \partial_{\mathbf{p}} Q^{\kappa}, \end{split}$$

where $\partial_p Q^{\kappa}$ denotes the parabolic boundary of the cylinder Q^{κ} for $\kappa = 0, 1, 2, \ldots$. Applying Proposition 3.3 to the equation of $u^{\kappa} - u$, we can obtain the the solvability and interior regularity of each u^{κ} .

Now, we *claim* that there is a constant $C = C(n, \lambda, \gamma)$ such that

(4.5)
$$[D^{m}(u^{\kappa} - u^{\kappa+1})]_{0;Q^{\kappa+2}} \le C\rho^{(2-m)\kappa-m}\varpi(\rho^{\kappa}), \quad m = 1, 2, \dots$$

To see this, we apply the second estimate in Proposition 3.1 (with f and g vanishing) to $u^{\kappa} - u^{\kappa+1}$ with $|\beta| = m, r = \rho^{\kappa+1}, \theta = 1/2$ and $p = \gamma$ to get

$$\|D^{m}(u^{\kappa} - u^{\kappa+1})\|_{0;Q^{\kappa+2}} \le C\rho^{-m\kappa-m} \left\| \oint_{Q^{\kappa+1}} (u^{\kappa} - u^{\kappa+1})^{2} \,\mathrm{d}X \right\|_{L^{\gamma/2}_{\omega}}^{1/2} =: I_{\kappa,m}$$

Here and in what follows, we denote $f_Q = \frac{1}{|Q|} \int_Q$ with |Q| being the Lebesgue measure of the set $Q \subset \mathbf{R}^{n+1}$.

On the other hand, it follows from Proposition 3.3 that

$$J_{\kappa} := \left\| \oint_{Q^{\kappa}} (u^{\kappa} - u)^2 \, \mathrm{d}X \right\|_{L^{\gamma/2}_{\omega}}^{1/2} \le C \rho^{2\kappa} \varpi(\rho^{\kappa}).$$

Combining the above we obtain

$$I_{\kappa,m} \le C\rho^{-m\kappa-m}(J_{\kappa} + J_{\kappa+1}) \le C\rho^{(2-m)\kappa-m}\varpi(\rho^{\kappa})$$

and thus the claim (4.5).

The estimate (4.5) with m = 2 gives (recalling $\rho = 1/2$)

$$\sum_{\kappa \ge 1} \| (u^{\kappa} - u^{\kappa+1})_{xx} \|_{0;Q^{\kappa+2}} \le C\rho^{-2} \sum_{\kappa \ge 1} \varpi(\rho^{\kappa}) \le 4C \int_0^1 \frac{\varpi(r)}{r} \,\mathrm{d}r < \infty,$$

which implies that $u_{xx}^{\kappa}(0)$ converges in L_{ω}^{γ} as $\kappa \to \infty$, (here $0 \in \mathbf{R}^{n+1}$). We shall prove that the limit is $u_{xx}(0)$. Since $\gamma \geq 2$, it suffices to show that

(4.6)
$$\lim_{\kappa \to \infty} \|u_{xx}^{\kappa}(0) - u_{xx}(0)\|_{L^{2}_{\omega}} = 0.$$

Applying the second estimate in Proposition 3.1 to $u^{\kappa} - u$ with m = n + 2, $|\beta| = 2, r = \rho^{\kappa}, \theta = 1/2$ and p = 2, we have

$$\begin{split} \sup_{Q^{\kappa+1}} \|u_{xx}^{\kappa} - u_{xx}\|_{L^{2}_{\omega}}^{2} &\leq C\rho^{-4\kappa} \mathbb{E} \oint_{Q^{\kappa}} |u^{\kappa} - u|^{2} \, \mathrm{d}X + C \mathbb{E} \oint_{Q^{\kappa}} |f(x,t) - f(0,t)|^{2} \, \mathrm{d}X \\ &+ C \mathbb{E} \oint_{Q^{\kappa}} \left(\|g(x,t) - g(0,t) - g_{x}(0,t) \cdot x\|^{2} + \|g_{x}(x,t) - g_{x}(0,t)\|^{2} \right) \, \mathrm{d}X \\ &+ C \sum_{k=1}^{n+1} \rho^{2\kappa k} \mathbb{E} \oint_{Q^{\kappa}} \left(|D^{k}f|^{2} + \|D^{k+1}g\|^{2} \right) \, \mathrm{d}X. \end{split}$$

According to the additional condition (A) on f and g, it is clear that the last three terms on the right-hand side tend to zero as $\kappa \to \infty$. Moreover, from Proposition 3.3 and (4.4) we have

$$\rho^{-4\kappa} \mathbb{E} \oint_{Q^{\kappa}} |u^{\kappa} - u|^2 dX$$

$$\leq C \mathbb{E} \oint_{Q^{\kappa}} \left(|f(x,t) - f(0,t)|^2 + \rho^{-2\kappa} |g(x,t) - g(0,t) - g_x(0,t) \cdot x|_{l^2}^2 \right) dX$$

$$\leq C \varpi (\rho^{\kappa})^2 \to 0, \quad \text{as } \kappa \to \infty.$$

Therefore, (4.6) is proved and $u_{xx}^{\kappa}(0)$ converges strongly to $u_{xx}(0)$ in L_{ω}^{γ} . Moreover, by means of (4.5), we have

(4.7)
$$\|u_{xx}^{\kappa}(0) - u_{xx}(0)\|_{L^{\gamma}_{\omega}} \leq \sum_{j \geq \kappa} \|(u^j - u^{j+1})_{xx}\|_{0;Q^{j+2}} \leq C \int_0^{\rho^{\kappa}} \frac{\varpi(r)}{r} \, \mathrm{d}r,$$

where $C = C(n, \lambda, \gamma)$.

Next we estimate the oscillation of u_{xx}^{κ} . Starting from $\kappa = 0$, u_{xx}^{0} satisfies the following homogeneous equation:

(4.8)
$$du_{xx}^{0} = a^{ij} D_{ij} u_{xx}^{0} dt + \sigma^{ik} D_{i} u_{xx}^{0} dw_{t}^{k} \quad \text{in } Q_{3/4}.$$

Using the second estimate in Proposition 3.1 (with f and g vanishing) to u_{xx}^0 , we have

$$\begin{aligned} \|D_x u_{xx}^0\|_{0;Q_{1/4}} + \|D_x^2 u_{xx}^0\|_{0;Q_{1/4}} &\leq C \|u_{xx}^0\|_{L_{\omega}^{\gamma} L_t^2 L_x^2(Q_{1/2})} \\ &\leq C (\|u_{xx}^0 - u_{xx}\|_{L_{\omega}^{\gamma} L_t^2 L_x^2(Q_{1/2})} + \|u_{xx}\|_{L_{\omega}^{\gamma} L_t^2 L_x^2(Q_{1/2})}). \end{aligned}$$

Then we apply the first estimate in Proposition 3.1 to u to get

$$\|u_{xx}\|_{L^{\gamma}_{\omega}L^{2}_{t}L^{2}_{x}(Q_{1/2})} \leq C(\|u\|_{L^{\gamma}_{\omega}L^{2}_{t}L^{2}_{x}(Q_{1})} + \|f\|_{L^{\gamma}_{\omega}L^{2}_{t}L^{2}_{x}(Q_{1})} + \|g\|_{L^{\gamma}_{\omega}L^{2}_{t}H^{1}_{x}(Q_{1})}),$$

and to $u^0 - u$ along with Proposition 3.3,

$$\begin{aligned} \|u_{xx}^{0} - u_{xx}\|_{L_{\omega}^{\gamma}L_{t}^{2}L_{x}^{2}(Q_{1/2})} &\leq C(\|u^{0} - u\|_{L_{\omega}^{\gamma}L_{t}^{2}L_{x}^{2}(Q_{1})} + \|f\|_{L_{\omega}^{\omega}L_{t}^{2}L_{x}^{2}(Q_{1})} + \|g\|_{L_{\omega}^{\omega}L_{t}^{2}H_{x}^{1}(Q_{1})}) \\ &\leq C(\|f\|_{L_{\omega}^{\omega}L_{t}^{2}L_{x}^{2}(Q_{1})} + \|g\|_{L_{\omega}^{\omega}L_{t}^{2}H_{x}^{1}(Q_{1})}). \end{aligned}$$

Therefore,

$$\begin{split} \|D_x u_{xx}^0\|_{0;Q_{1/4}} + \|D_x^2 u_{xx}^0\|_{0;Q_{1/4}} \\ &\leq C(\|u\|_{L_{\omega}^{\gamma}L_t^2 L_x^2(Q_1)} + \|f\|_{L_{\omega}^{\gamma}L_t^2 L_x^2(Q_1)} + \|g\|_{L_{\omega}^{\gamma}L_t^2 H_x^1(Q_1)}) \leq CM_1. \end{split}$$

Hence, for $-1/16 < s \leq t \leq 0$ and $x \in B_{1/4}$,

$$\begin{aligned} \|u_{xx}^{0}(x,t) - u_{xx}^{0}(x,s)\|_{L_{\omega}^{\gamma}} &= \left\| \int_{s}^{t} a^{ij} D_{ij} u_{xx}^{0} \,\mathrm{d}\tau + \int_{s}^{t} \sigma^{ik} D_{i} u_{xx}^{0} \,\mathrm{d}w_{\tau}^{k} \right\|_{L_{\omega}^{\gamma}} \\ &\leq C\sqrt{t-s} (\|Du_{xx}^{0}\|_{0;Q_{1/4}} + \|D^{2} u_{xx}^{0}\|_{0;Q_{1/4}}) \leq C\sqrt{t-s} M_{1}, \end{aligned}$$

where $C = C(n, \lambda, \gamma)$. Thus, we obtain

(4.9)
$$||u_{xx}^{0}(X) - u_{xx}^{0}(Y)||_{L_{\omega}^{\gamma}} \le CM_{1}|X - Y|_{\mathbf{p}}, \quad \forall X, Y \in Q_{1/4}.$$

To deal with u_{xx}^{κ} with $\kappa \geq 1$, we denote

$$h^{\iota} = u^{\iota} - u^{\iota-1}, \quad \text{for } \iota = 1, 2, \dots, \kappa.$$

Then h^{ι} satisfies

(4.10)
$$\mathrm{d}h^{\iota} = a^{ij}h^{\iota}_{ij}\,\mathrm{d}t + \sigma^{ik}h^{\iota}_{i}\,\mathrm{d}w^{k}_{t} \quad \text{in } Q^{\iota}$$

By (4.5) we have

$$\rho^{-\iota} [\![D^3 h^{\iota}]\!]_{0;Q^{\iota+1}} + [\![D^4 h^{\iota}]\!]_{0;Q^{\iota+1}} \le C \rho^{-2\iota} \varpi(\rho^{\iota-1})$$

Hence, for $-\rho^{2(\kappa+1)} \le t \le 0$ and $|x| \le \rho^{\kappa+1}$,

$$\|h_{xx}^{\iota}(x,0)-h_{xx}^{\iota}(0,0)\|_{L_{\omega}^{\gamma}}\leq C\rho^{\kappa-\iota}\varpi(\rho^{\iota-1})$$

and

$$\begin{aligned} \|h_{xx}^{\iota}(x,t) - h_{xx}^{\iota}(x,0)\|_{L_{\omega}^{\gamma}} &= \left\|\int_{t}^{0} a^{ij} D_{ij} h_{xx}^{\iota} \,\mathrm{d}\tau + \int_{t}^{0} \sigma^{ik} D_{i} h_{xx}^{\iota} \,\mathrm{d}w_{\tau}^{k}\right\|_{L_{\omega}^{\gamma}} \\ &\leq C \rho^{2\kappa} \|D^{4} h^{\iota}\|_{0;Q^{\iota+1}} + C \rho^{\kappa} \|D^{3} h^{\iota}\|_{0;Q^{\iota+1}} \leq C \rho^{\kappa-\iota} \varpi(\rho^{\iota-1}). \end{aligned}$$

Let $Y = (y, s) \in Q_{1/4}$, and $\tilde{\kappa} \in \mathbf{N}$ such that

$$\delta := |Y|_{\mathbf{p}} \in [\rho^{\tilde{\kappa}+2}, \rho^{\tilde{\kappa}+1}).$$

Combining the last two estimates and (4.9), we can obtain $\| \tilde{k}_{k}(\mathbf{V}) - \tilde{k}_{k}(\mathbf{O}) \|_{\mathbf{V}} \le \| \tilde{k}^{-1}(\mathbf{V}) - \tilde{k}^{-1}(\mathbf{O}) \|_{\mathbf{V}}$

$$\begin{split} \|u_{xx}^{\kappa}(Y) - u_{xx}^{\kappa}(0)\|_{L_{\omega}^{\gamma}} &\leq \|u_{xx}^{\kappa-1}(Y) - u_{xx}^{\kappa-1}(0)\|_{L_{\omega}^{\gamma}} + \|h_{xx}^{\kappa}(Y) - h_{xx}^{\kappa}(0)\|_{L_{\omega}^{\gamma}} \\ &\leq \|u_{xx}^{0}(Y) - u_{xx}^{0}(0)\|_{L_{\omega}^{\gamma}} + \sum_{\iota=1}^{\tilde{\kappa}} \|h_{xx}^{\iota}(Y) - h_{xx}^{\iota}(0)\|_{L_{\omega}^{\gamma}} \\ &\leq CM_{1}\rho^{\tilde{\kappa}+1} + C\sum_{\iota=1}^{\tilde{\kappa}} \rho^{\tilde{\kappa}-\iota}\varpi(\rho^{\iota-1}) \\ &\leq CM_{1}\rho^{\tilde{\kappa}+2} + C\rho^{\tilde{\kappa}+2}\int_{\rho^{\tilde{\kappa}}}^{1}\frac{\varpi(r)}{r^{2}}\,\mathrm{d}r \\ &\leq C\delta M_{1} + C\delta\int_{\delta}^{1}\frac{\varpi(r)}{r^{2}}\,\mathrm{d}r. \end{split}$$

By virtue of (4.7) we have the following decomposition

$$(4.11) \|u_{xx}(Y) - u_{xx}(0)\|_{L_{\omega}^{\gamma}}
\leq \|u_{xx}^{\tilde{\kappa}}(Y) - u_{xx}^{\tilde{\kappa}}(0)\|_{L_{\omega}^{\gamma}} + \|u_{xx}^{\tilde{\kappa}}(0) - u_{xx}(0)\|_{L_{\omega}^{\gamma}} + \|u_{xx}^{\tilde{\kappa}}(Y) - u_{xx}(Y)\|_{L_{\omega}^{\gamma}}
\leq C \left[\delta M_{1} + \int_{0}^{4\delta} \frac{\varpi(r)}{r} \, \mathrm{d}r + \delta \int_{\delta}^{1} \frac{\varpi(r)}{r^{2}} \, \mathrm{d}r\right] + \|u_{xx}^{\tilde{\kappa}}(Y) - u_{xx}(Y)\|_{L_{\omega}^{\gamma}}.$$

It remains to estimate the last term in the above inequality. To this end, we consider the sequence of equations

$$du^{Y,\kappa} = [a^{ij}u^{Y,\kappa}_{ij} + f(y,t)] dt + [\sigma^{ik}u^{Y,\kappa}_i + g^k(y,t) + g^k_x(y,t) \cdot x] dw^k_t \quad \text{in } Q^{\kappa}(Y),$$

$$u^{Y,\kappa} = u \quad \text{on } \partial_{\mathbf{p}}Q^{\kappa}(Y) \quad \text{with } \kappa = 0, 1, \dots, \tilde{\kappa} - 1, \tilde{\kappa} + 2, \dots;$$

the equations associated with $\tilde{\kappa}$ and $\tilde{\kappa}+1$ are replaced by the following single equation

$$\begin{aligned} \mathrm{d}u^{Y,\tilde{\kappa}} &= \left[a^{ij}u^{Y,\tilde{\kappa}}_{ij} + f(y,t)\right]\mathrm{d}t + \left[\sigma^{ik}u^{Y,\tilde{\kappa}}_{i} + g^{k}(y,t) + g^{k}_{x}(y,t) \cdot x\right]\mathrm{d}w^{k}_{t} \quad \text{in } Q^{\tilde{\kappa}}(0), \\ u^{Y,\tilde{\kappa}} &= u \quad \text{on } \partial_{\mathbf{p}}Q^{\tilde{\kappa}}(0). \end{aligned}$$

As $|Y|_{\mathbf{p}} \in [\rho^{\tilde{\kappa}+2}, \rho^{\tilde{\kappa}+1})$, it is easily seen that $Q^{\tilde{\kappa}+2}(Y) \subset Q^{\tilde{\kappa}}(0) \subset Q^{\tilde{\kappa}-1}(Y)$. So analogously to proving (4.7) but only with minor changes, one can derive

(4.12)
$$\|u_{xx}^{Y,\tilde{\kappa}}(Y) - u_{xx}(Y)\|_{L^{\gamma}_{\omega}} \le C \int_{0}^{\rho^{\kappa}} \frac{\overline{\omega}(r)}{r} \,\mathrm{d}r,$$

where $C = C(n, \lambda, \gamma)$. On the other hand, applying Proposition 3.1 to the equation satisfied by $u^{Y,\tilde{\kappa}} - u^{\tilde{\kappa}}$, and using (4.4), we have

$$\|u_{xx}^{Y,\tilde{\kappa}}(Y) - u_{xx}^{\tilde{\kappa}}(Y)\|_{L^{\gamma}_{\omega}} \le C\big(\|u^{Y,\tilde{\kappa}} - u^{\tilde{\kappa}}\|_{L^{\gamma}_{\omega}L^{2}_{t}L^{2}_{x}(Q^{\tilde{\kappa}})} + \varpi(\delta)\big),$$

while by Proposition 3.3,

$$\|u^{Y,\tilde{\kappa}} - u^{\tilde{\kappa}}\|_{L^{\gamma}_{\omega}L^2_t L^2_x(Q_1)} \le C\varpi(\delta).$$

Thus

$$\|u_{xx}^{\tilde{\kappa}}(Y) - u_{xx}(Y)\|_{L^{\gamma}_{\omega}} \le C\varpi(\delta) + C \int_{0}^{4\delta} \frac{\varpi(r)}{r} \,\mathrm{d}r$$

Substituting the above estimate into (4.11), we then complete the proof.

Remark 4.1. Consider the model equation of divergence-form

(4.13)
$$du = (a^{ij}u_j + f^i)_i dt + (\sigma^{ik}u_i + g^k) dw_t^k.$$

With the help of the following approximation sequence

$$\begin{split} \mathrm{d} u^{\kappa} &= a^{ij} u_{ij}^{\kappa} \,\mathrm{d} t + \left[\sigma^{ik} u_{i}^{\kappa} + g^{k}(0,t)\right] \mathrm{d} w_{t}^{k} \quad \text{in } Q^{\kappa}, \\ u^{\kappa} &= u \quad \text{on } \partial_{\mathbf{p}} Q^{\kappa}, \end{split}$$

we can similarly obtain an interesting estimate

$$\llbracket u \rrbracket_{(1+\alpha,\alpha/2);Q_{1/4}} \le C \bigl(\llbracket u \rrbracket_{0;Q_1} + \sum_i \llbracket f^i \rrbracket_{\alpha;Q_1} + \llbracket g \rrbracket_{\alpha;Q_1} \bigr),$$

provided the right-hand side is finite. The result on the model equation (4.13) can help us establish a $C^{1+\alpha}$ estimate for more general equations, which will be discussed in a separate work.

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5. GLOBAL HÖLDER ESTIMATES FOR GENERAL EQUATIONS

This section is devoted to the proof of Theorem 1.3. First we state two technical lemmas.

Lemma 5.1. Let φ be a bounded nonnegative function defined on [0,T] satisfying

(5.1)
$$\varphi(t) \le \theta \varphi(s) + \sum_{i=1}^{m} A_i (s-t)^{-\delta_i}, \quad \forall \, 0 \le t < s \le T$$

for some nonnegative constants θ , δ_i and A_i (i = 1, ..., m), where $\theta < 1$. Then

$$\varphi(0) \le C \sum_{i=1}^{m} A_i T^{-\delta_i},$$

where C depends only on $\delta_1, \ldots, \delta_m$ and θ .

Proof. We may suppose T = 1, otherwise let $\tilde{\varphi}(t) = \varphi(Tt)$. Then (5.1) implies

$$\varphi(t) \le \theta \varphi(s) + A(s-t)^{-\delta}, \quad \forall 0 \le t < s \le 1,$$

where $\delta := \max_{1 \le i \le m} \delta_i$ and $A := A_1 + \cdots + A_m$. It suffices to consider $\delta > 0$. Take $\tau \in (0, 1)$ such that $\varepsilon := \theta \tau^{-\delta} < 1$, and set $t_0 = 0$, $t_{j+1} = t_j + (1 - \tau)\tau^j$ for $j = 0, 1, \ldots$ Then

$$\theta^{j}\varphi(t_{j}) \leq \theta^{j+1}\varphi(t_{j+1}) + A(1-\tau)^{-\delta}(\theta\tau^{-\delta})^{j} = \theta^{j+1}\varphi(t_{j+1}) + \varepsilon^{j}A(1-\tau)^{-\delta}$$

By iteration, we gain

$$\varphi(0) \le \theta^k \varphi(t_k) + (1 + \varepsilon + \dots + \varepsilon^{k-1}) A (1 - \tau)^{-\delta}$$

$$\le \theta^k \varphi(t_k) + (1 - \varepsilon)^{-1} A (1 - \tau)^{-\delta}.$$

By letting $k \to \infty$, we conclude the proof.

Lemma 5.2. Let $B_R = \{x \in \mathbb{R}^n : |x| < R\}$ with R > 0, E be a Banach space, $p \ge 1$, and $0 \le s < r$. There exists a positive constant C, depending only on n and p, such that

(5.2)
$$[u]_{s;B_R}^E \le C\varepsilon^{r-s} [u]_{r;B_R}^E + C\varepsilon^{-s-n/p} ||u||_{L^p(B_R;E)}$$

for any $u \in C^r(B_R; E)$ and $\varepsilon \in (0, R)$.

Proof. Let us consider R = 1 first. Following the proof of classical interpolation inequalities for Hölder norms (see [GT01, Lemma 6.35] or [Kry96a, Theorem 3.2.1]), one can derive

(5.3)
$$[u]_{s;B_1}^E \le C\varepsilon^{r-s} [u]_{r;B_1}^E + C\varepsilon^{-s} |u|_{0;B_1}^E .$$

Consider the case of $r \leq 1$ first. For arbitrary $x \in B_1$, we select a ball $B_{\varepsilon}(y) = \{z : |y - z| < \varepsilon\} \subset B_1$ such that $x \in B_{\varepsilon}(y)$. Then we compute

$$\begin{aligned} \|u(x)\|_{E} &= \int_{B_{\varepsilon}(y)} \|u(x)\|_{E} \,\mathrm{d}z \leq \int_{B_{\varepsilon}(y)} \|u(x) - u(z)\|_{E} \,\mathrm{d}z + \int_{B_{\varepsilon}(y)} \|u(z)\|_{E} \,\mathrm{d}z \\ &\leq C\varepsilon^{r} [u]_{r;B_{1}}^{E} + C\varepsilon^{-n/p} \|u\|_{L^{p}(B_{1};E)}, \end{aligned}$$

which yields (5) with s = 0 and $r \le 1$. For the case of r > 1, we have

$$|u|_{0;B_{1}}^{E} \leq C\varepsilon[u]_{1;B_{1}}^{E} + C\varepsilon^{-n/p} ||u||_{L^{p}(B_{1};E)}$$

$$\leq C\varepsilon\varepsilon_{1}^{r-1}[u]_{r;B_{1}}^{E} + C\varepsilon\varepsilon_{1}^{-1}|u|_{0;B_{1}}^{E} + C\varepsilon^{-n/p} ||u||_{L^{p}(B_{1};E)}$$

Choosing $\varepsilon_1 = 2C\varepsilon$, we get (5) with s = 0 and r > 1. Finally, for the general case, we derive

(5.4)
$$|u|_{s;B_1}^E \leq C\varepsilon^{r-s}[u]_{r;B_1}^E + C\varepsilon^{-s}|u|_{0;B_1}^E \leq C\varepsilon^{r-s}[u]_{r;B_1}^E + C\varepsilon^{-s-n/p}||u||_{L^p(B_1;E)}.$$

The case of R = 1 is proved.

Now we turn to the general R > 0. With v(x) := u(Rx) we have that $[v]_{r;B_1}^E = R^r[u]_{r;B_R}^E$ and $\|v\|_{L^p(B_1;E)} = R^{-n/p} \|u\|_{L^p(B_R;E)}$. Applying (5.4) to v we can obtain that

$$[u]_{s;B_R}^E \le C(R\varepsilon)^{r-s} [u]_{r;B_R}^E + C(R\varepsilon)^{-s-n/p} ||u||_{L^p(B_R;E)}$$

for any $u \in C^r(B_R; E)$ and $\varepsilon \in (0, 1)$. The proof is complete.

We are now ready to prove Theorem 1.3.

Proof of Theorem 1.3. Let $\rho/2 \leq r < R \leq \rho$ with $\rho \in (0, 1/4)$ to be specified. Take a nonnegative function $\zeta \in C_0^{\infty}(\mathbf{R}^n)$ such that $\zeta(x) = 1$ when $|x| \leq r$; $\zeta(x) = 0$ when |x| > R, and for $\delta \geq 0$,

$$[\zeta]_{\delta;\mathbf{R}^n} \le C(R-r)^{-\delta}.$$

Set $v = \zeta u$, and

$$a^{ij}_{\circ}(t) = a^{ij}(0,t), \quad \sigma^{ik}_{\circ}(t) = \sigma^{ik}(0,t).$$

Then v satisfies

(5.5)
$$\mathrm{d}v = (a_{\circ}^{ij}v_{ij} + \tilde{f})\,\mathrm{d}t + (\sigma_{\circ}^{ik}v_i + \tilde{g}^k)\,\mathrm{d}w_t^k,$$

where

$$\tilde{f} = (a^{ij} - a^{ij}_{\circ})\zeta u_{ij} + (b^i\zeta - 2a^{ij}\zeta_j)u_i + (c\zeta - a^{ij}\zeta_{ij} - b^i\zeta_i)u + \zeta f,$$

$$\tilde{g}^k = (\sigma^{ik} - \sigma^{ik}_{\circ})\zeta u_i + (\nu^k\zeta - \sigma^{ik}\zeta_i)u + \zeta g^k.$$

For a positive number τ , we set $\mathcal{Q}_{R,\tau} = B_R \times (0,\tau)$ and define

$$M_{x,r}^{\tau}(u) = \sup_{0 \le t \le \tau} \left(\oint_{B_r(x)} \mathbb{E} |u(t,y)|^{\gamma} \, \mathrm{d}y \right)^{1/\gamma}, \quad M_r^{\tau}(u) = \sup_{x \in \mathbf{R}^n} M_{x,r}^{\tau}(u).$$

Then by Lemma 5.2,

$$\begin{split} \|\tilde{f}\|_{\alpha;\mathcal{Q}_{R,\tau}} &\leq (\varepsilon + KR^{\alpha}) [\![u]\!]_{2+\alpha;\mathcal{Q}_{R,\tau}} + C(R-r)^{-2-\alpha-n/\gamma} M_{0,R}^{\tau}(u) \\ &+ [\![f]\!]_{\alpha;\mathcal{Q}_{R,\tau}} + C(R-r)^{-\alpha} [\![f]\!]_{0;\mathcal{Q}_{R,\tau}}, \\ \|\tilde{g}\|_{1+\alpha;\mathcal{Q}_{R,\tau}} &\leq (\varepsilon + KR^{\alpha}) [\![u]\!]_{2+\alpha;\mathcal{Q}_{R,\tau}} + C(R-r)^{-2-\alpha-n/\gamma} M_{0,R}^{\tau}(u) \\ &+ [\![g]\!]_{1+\alpha;\mathcal{Q}_{R,\tau}} + C(R-r)^{-1-\alpha} [\![g]\!]_{0;\mathcal{Q}_{R,\tau}}. \end{split}$$

where $C = C(n, K, \gamma, \varepsilon, \rho)$. Take $\rho, \varepsilon > 0$ so small that $\varepsilon + K\rho^{\alpha} \leq 1/4$, then by virtue of Corollary 4.2, we obtain that for any $\rho/2 \leq r < R \leq \rho$,

$$\begin{split} \llbracket u \rrbracket_{(2+\alpha,\alpha/2);\mathcal{Q}_{r,\tau}} &\leq \frac{1}{2} \llbracket u \rrbracket_{2+\alpha;\mathcal{Q}_{R,\tau}} + C(R-r)^{-2-\alpha-n/\gamma} M_{0,R}^{\tau}(u) + \llbracket f \rrbracket_{\alpha;\mathcal{Q}_{R,\tau}} \\ &+ \llbracket g \rrbracket_{1+\alpha;\mathcal{Q}_{R,\tau}} + C(R-r)^{-\alpha} \llbracket f \rrbracket_{0;\mathcal{Q}_{R,\tau}} + C(R-r)^{-1-\alpha} \llbracket g \rrbracket_{0;\mathcal{Q}_{R,\tau}}. \end{split}$$

By Lemma 5.1, we gain

(5.6)
$$\llbracket u \rrbracket_{(2+\alpha,\alpha/2);\mathcal{Q}_{\rho/2,\tau}} \le C \Big(M_{0,\rho}^{\tau}(u) + \llbracket f \rrbracket_{\alpha;\mathcal{Q}_{\rho,\tau}} + \llbracket g \rrbracket_{1+\alpha;\mathcal{Q}_{\rho,\tau}} \Big).$$

Similarly, with $\mathcal{Q}_{\rho/2,\tau}(x) := B_{\rho/2}(x) \times (0,\tau)$ for any point $x \in \mathbf{R}^n$, we have

$$\sup_{x \in \mathbf{R}^n} \llbracket u \rrbracket_{(2+\alpha,\alpha/2);\mathcal{Q}_{\rho/2,\tau}(x)} \le C \Big(M_{\rho}^{\tau}(u) + \llbracket f \rrbracket_{\alpha;\mathcal{Q}_{\tau}} + \llbracket g \rrbracket_{1+\alpha;\mathcal{Q}_{\tau}} \Big),$$

which along with the localization property of Hölder norms (cf. [Kry96a, Lemma 4.1.1]) and Lemma 5.2, we have

(5.7)
$$\llbracket u \rrbracket_{(2+\alpha,\alpha/2);\mathcal{Q}_{\tau}} \leq C \sup_{x \in \mathbf{R}^{n}} \left(\llbracket u \rrbracket_{(2+\alpha,\alpha/2);\mathcal{Q}_{\rho/2,\tau}(x)} + \llbracket u \rrbracket_{0;\mathcal{Q}_{\rho/2,\tau}(x)} \right)$$
$$\leq C \left(\sup_{x \in \mathbf{R}^{n}} \llbracket u \rrbracket_{(2+\alpha,\alpha/2);\mathcal{Q}_{\rho/2,\tau}(x)} + M_{\rho/2}^{\tau}(u) \right)$$
$$\leq C \left(M_{\rho}^{\tau}(u) + \llbracket f \rrbracket_{\alpha;\mathcal{Q}_{\tau}} + \llbracket g \rrbracket_{1+\alpha;\mathcal{Q}_{\tau}} \right),$$

where $C = C(n, \lambda, \gamma, \alpha)$ is a generic constant.

To estimate $M^{\tau}_{\rho}(u)$, we apply Itô's formula to compute

$$\begin{aligned} \mathrm{d}|u|^{\gamma} &= \gamma |u|^{\gamma-2} u(a^{ij}u_{ij} + b^i u_i + cu + f) \,\mathrm{d}t \\ &+ \frac{\gamma(\gamma-1)}{2} |u|^{\gamma-2} |\sigma^{ik}u_i + \nu^k u + g^k|^2 \,\mathrm{d}t + \mathrm{d}m_t \end{aligned}$$

where m_t is a martingale. Integrating in $\mathcal{Q}_{\rho,\tau} \times \Omega$, and using the Hölder and Sobolev–Gagliargo–Nirenberg inequalities, we get

$$\mathbb{E}\int_{B_{\rho}}|u(x,\tau)|^{\gamma}\,\mathrm{d}x \leq C_{1}\mathbb{E}\int_{\mathcal{Q}_{\rho,\tau}}\left(|u_{xx}|^{\gamma}+|u|^{\gamma}+|f|^{\gamma}+|g|^{\gamma}\right)\mathrm{d}x\mathrm{d}t.$$

Thus,

 $M_{0,\rho}^{\tau}(u) \le C_1 \tau([]u[]_{2;\mathcal{Q}_{\rho,\tau}} + []f]_{0;\mathcal{Q}_{\tau}} + []g[]_{0;\mathcal{Q}_{\tau}}),$

where $C_1 = C_1(n, \lambda, \gamma)$. Letting $\tau = (2CC_1)^{-1}$, the above inequality along with (5.7) yields

(5.8)
$$\|u\|_{(2+\alpha,\alpha/2);\mathcal{Q}_{\tau}} \leq C_0 (\|f\|_{\alpha;\mathcal{Q}_{\tau}} + \|g\|_{1+\alpha;\mathcal{Q}_{\tau}}),$$

where $C_0 = C_0(n, \lambda, \gamma, \alpha)$.

Let us conclude the proof by induction. For S > 0, assume that there is a constant C_S such that

(5.9)
$$[\![u]\!]_{(2+\alpha,\alpha/2);\mathcal{Q}_S} \leq C_S([\![f]\!]_{\alpha;\mathcal{Q}_S} + [\![g]\!]_{1+\alpha;\mathcal{Q}_S}).$$

With $u^{S} := u(\cdot, S)$, it is easily seen that $v := (u - u^{S})$ satisfies $dv = [a^{ij}v_{ij} + b^{i}v_{i} + cv + (f + a^{ij}u_{ij}^{S} + b^{i}u_{i}^{S} + cu^{S})] dt$ $+ [\sigma^{ik}v_{i} + \nu^{k}v + (g^{k} + \sigma^{ik}u_{i}^{S} + \nu^{k}u^{S})] dw_{t}^{k}, \text{ on } \mathbf{R}^{n} \times (S, \infty),$ $v(S, x) = 0, \qquad x \in \mathbf{R}^{n}.$

Applying (5.8) to this equation and with (5.9) in mind, we have

$$\begin{aligned} \|u\|_{(2+\alpha,\alpha/2);\mathcal{Q}_{S+\tau}} &\leq \|v\|_{(2+\alpha,\alpha/2);\mathcal{Q}_{S+\tau}} + \|u\|_{(2+\alpha,\alpha/2);\mathcal{Q}_{S}} \\ &\leq C_0 \big(\|f\|_{\alpha;\mathcal{Q}_{S+\tau}} + \|g\|_{1+\alpha;\mathcal{Q}_{S+\tau}}\big) + C_0 N \|u\|_{(2+\alpha,\alpha/2);\mathcal{Q}_{S}} \\ &\leq C_0 (1 + NC_S) \big(\|f\|_{\alpha;\mathcal{Q}_{S+\tau}} + \|g\|_{1+\alpha;\mathcal{Q}_{S+\tau}}\big), \end{aligned}$$

where N is a constant depending only on n, λ, γ and K. Hence,

$$C_{S+\tau} \le C_0 (1 + NC_S).$$

As τ is fixed, by iteration we have $C_S \leq Ce^{CS}$, where $C = C(n, \lambda, \gamma, \alpha, K)$. This concludes the proof of estimate (2.2) and thus Theorem 1.3.

Appendix

In this section we prove some properties and approximation of L^{γ}_{ω} -valued continuous and differentiable functions. Let \mathcal{O} be a simply connected domain in \mathbb{R}^n . Denote $C(\mathcal{O}; L^{\gamma}_{\omega})$ the set of all L^{γ}_{ω} -valued strongly continuous functions defined on \mathcal{O} such that $\sup_{x \in \mathcal{O}} \mathbb{E}[|u(x)|^{\gamma}] < \infty$, and $C^m(\mathcal{O}; L^{\gamma}_{\omega})$ the set of all $C(\mathcal{O}; L^{\gamma}_{\omega})$ functions whose strong derivatives up to order m all exit and belong to $C(\mathcal{O}; L^{\gamma}_{\omega})$, where m is a nonnegative integer. In view of a known result (see [DPZ92, Proposition 3.6]), every function in $C(\mathcal{O}; L^{\gamma}_{\omega})$ has a modification jointly measurable with respect to $x \in \mathcal{O}$ and $\omega \in \Omega$; we will always choose this modification.

In what follows, we denote Du to the strong derivatives of an L^{γ}_{ω} -valued differentiable function u, and ∂u to be the classical derivatives if exist.

Lemma A.1. If $u \in L^{\gamma}(\Omega; C^m(\mathcal{O}))$, then $u \in C^m(\mathcal{O}; L^{\gamma}_{\omega})$, and $D^{\beta}u = \partial^{\beta}u$ for any multi-index β with $|\beta| \leq m$.

Proof. It follows from the dominated convergence theorem that $L^{\gamma}(\Omega; C(\mathcal{O})) \subset C(\mathcal{O}; L^{\gamma}_{\omega})$. For $u \in L^{\gamma}(\Omega; C^{1}(\mathcal{O}))$, we know that $\partial u \in C(\mathcal{O}; L^{\gamma}_{\omega})$, and by Jensen's inequality and Fubini's theorem,

$$\mathbb{E}|r^{-1}[u(x+re_i)-u(x)] - \partial_i u(x)|^{\gamma} = \mathbb{E}\left|\int_0^1 \left[\partial_i u(x+sre_i) - \partial_i u(x)\right] \mathrm{d}s\right|^{\gamma}$$
$$\leq \int_0^1 \mathbb{E}\left|\partial_i u(x+sre_i) - \partial_i u(x)\right|^{\gamma} \mathrm{d}s \to 0 \quad \text{as } r \to 0.$$

Thus, $Du = \partial u$ and $u \in C^1(\mathcal{O}; L^{\gamma}_{\omega})$. The lemma is concluded by induction.

Let $\varphi = \tilde{\varphi} / \int_{\mathbf{R}^n} \tilde{\varphi}$ with $\tilde{\varphi}(x) := e^{-1/(1-|x|^2)} \mathbf{1}_{B_1}(x)$ for $x \in \mathbf{R}^n$. Define $\varphi^{\varepsilon} = \varepsilon^{-n} \varphi(x/\varepsilon)$ with $\varepsilon > 0$. It is easily seen that

$$|D^m \varphi^{\varepsilon}| \le C \varepsilon^{-n-m}, \quad \int_{\mathbf{R}^n} |D^m \varphi^{\varepsilon}|^{\gamma} \le C \varepsilon^{-(\gamma-1)n-\gamma m}$$

for all $m \in \mathbf{N}$ and $\gamma \geq 2$. Now we mollify a function $u \in C(\mathbf{R}^n; L^{\gamma}_{\omega})$ using φ^{ε} :

(A.1)
$$u^{\varepsilon} := \varphi^{\varepsilon} * u = \int \varphi^{\varepsilon} (\cdot - y) u(y) \, \mathrm{d}y.$$

It is easily seen from Fubini's theorem that $u^{\varepsilon} \in C(\mathbf{R}^n; L^{\gamma}_{\omega})$. In view of Theorem 34.B in [Hal74], $u(x, \omega)$ is a measurable function in x for each ω , and $u(\cdot, \omega) \in L^{\gamma}_{\text{loc}}(\mathbf{R}^n)$ for almost every ω . Thus, by the property of mollifiers, $u^{\varepsilon}(\cdot, \omega) \in C^{\infty}(\mathbf{R}^n)$ for almost every ω .

Lemma A.2. If $u \in C(\mathbf{R}^n; L^{\gamma}_{\omega})$, then $u^{\varepsilon} \in \bigcap_{m \in \mathbf{N}} C^m(\mathbf{R}^n; L^{\gamma}_{\omega})$, and it restricted on any B_R belongs to $L^{\gamma}(\Omega; C^k(B_R))$ for all $k \in \mathbf{N}$.

Proof. By Hölder's inequality and Fubini's theorem we have

$$\mathbb{E}|\partial^{m}u^{\varepsilon}(x)|^{\gamma} = \mathbb{E}\left|\int_{B_{\varepsilon}(x)}\partial^{m}\varphi^{\varepsilon}(x-y)u(y)\,\mathrm{d}y\right|^{\gamma}$$

$$\leq |B_{\varepsilon}|^{\gamma-1}\int_{B_{\varepsilon}}|\partial^{m}\varphi^{\varepsilon}(y)|^{\gamma}\,\mathrm{d}y\Big(\sup_{x\in\mathbf{R}^{n}}\mathbb{E}|u(x)|^{\gamma}\Big) \leq C\varepsilon^{-p(m+1)}[]u[]_{0;\mathbf{R}^{n}}^{\gamma},$$

which implies that $u^{\varepsilon} \in \bigcap_{m \in \mathbb{N}} C^m(\mathbb{R}^n; L^{\gamma}_{\omega})$ and also that $u^{\varepsilon} \in L^{\gamma}(\Omega; W^{k,\gamma}_{\text{loc}}(\mathbb{R}^n))$ for any $k \in \mathbb{N}$. By Sobolev's embedding theorem, $u^{\varepsilon} \in L^{\gamma}(\Omega; C^k(B_R))$ for any $k \in \mathbb{N}$ and R > 0.

Remark A.1. If u also depends on time, say $u \in C(\mathbf{R}^n \times \mathbf{R}; L^{\gamma}_{\omega})$, then the above lemma implies that $u^{\varepsilon}(\cdot, t) \in L^{\gamma}(\Omega; C^k(B_R))$ and $\sup_{t \in [-R,R]} \mathbb{E}|u^{\varepsilon}|_{k;B_R}^{\gamma} < \infty$ for all $k \in \mathbf{N}$ and R > 0. In particular, $u^{\varepsilon} \in L^{\gamma}_{\omega}L^2_t H^k_x(Q_R) \cap C^k_x(Q_R; L^{\gamma}_{\omega})$ for all $k \in \mathbf{N}$ and R > 0.

Lemma A.3. If $u \in C(\mathbf{R}^n; L^{\gamma}_{\omega})$, then

(A.2)
$$\lim_{\varepsilon \to 0} \mathbb{E}[|u^{\varepsilon}(x) - u(x)|^{\gamma}] = 0 \quad \forall x \in \mathbf{R}^{n};$$

if, in addition, u is uniformly strongly continuous, namely

$$\lim_{\varepsilon \to 0} \sup_{x \in \mathbf{R}^n} \max_{|y| \le \varepsilon} \mathbb{E}[|u(x+y) - u(x)|^{\gamma}] = 0,$$

then the convergence (A.2) is uniform with respect to $x \in \mathbf{R}^n$.

Proof. Using the continuity of u we have

$$\mathbb{E}|u^{\varepsilon}(x) - u(x)|^{\gamma} = \mathbb{E}\left|\int_{B_{\varepsilon}(x)} \varphi^{\varepsilon}(x - y)[u(y) - u(x)] \,\mathrm{d}y\right|^{\gamma}$$

$$\leq |B_{\varepsilon}|^{\gamma - 1} \int_{B_{\varepsilon}} |\varphi^{\varepsilon}(y)|^{\gamma} \,\mathrm{d}y \left(\max_{|y| \leq \varepsilon} \mathbb{E}|u(x + y) - u(x)|^{\gamma}\right)$$

$$\leq C \max_{|y| \leq \varepsilon} \mathbb{E}|u(x + y) - u(x)|^{\gamma} \to 0 \quad \text{as } \varepsilon \to 0.$$

Then the lemma is easily concluded.

Lemma A.4. If $u \in C^1(\mathbf{R}^n; L^{\gamma}_{\omega})$, then $u^{\varepsilon} \in C^1(\mathbf{R}^n; L^{\gamma}_{\omega})$ and $Du^{\varepsilon} = \varphi^{\varepsilon} * Du$.

Proof. We compute that

$$\mathbb{E} \left| r^{-1} [u^{\varepsilon}(x+re_{i})-u^{\varepsilon}(x)] - \varphi^{\varepsilon} * D_{i}u(x) \right|^{\gamma}$$

$$= \mathbb{E} \left| \int_{B_{\varepsilon}(x)} \varphi^{\varepsilon}(x-y) \left(\frac{u(y+re_{i})-u(y)}{r} - D_{i}u(y) \right) dy \right|^{\gamma}$$

$$\leq \int_{B_{\varepsilon}(x)} \int_{0}^{1} \varphi^{\varepsilon}(x-y) \mathbb{E} \left| D_{i}u(y+sre_{i}) - D_{i}u(y) \right|^{\gamma} dsdy.$$

Since $\lim_{r\to 0} \mathbb{E} |D_i u(y + sre_i) - D_i u(y)|^{\gamma} = 0$ and $\mathbb{E} |D_i u(y + sre_i) - D_i u(y)|^{\gamma} \leq 2^{\gamma} ||u||_{1:\mathbf{R}^n}^{\gamma}$, the lemma is concluded by the dominated convergence theorem. \Box

The following consequence is straightforward.

Corollary A.5. If $u \in C^m(\mathbf{R}^n; L^{\gamma}_{\omega})$, then $D^m u^{\varepsilon}(x)$ converges to $D^m u(x)$ in L^{γ}_{ω} for each $x \in \mathbf{R}^n$, as ε tends to zero; if, in addition, $D^m u$ is uniformly strongly continuous, then the convergence is uniform with respect to $x \in \mathbf{R}^n$.

The final lemma concerns the convergence of Hölder norms.

Lemma A.6. If $u \in C^{\alpha}(\mathbf{R}^n; L^{\gamma}_{\omega})$ with $\alpha \in (0, 1)$, then $u^{\varepsilon} \in C^{\alpha}(\mathbf{R}^n; L^{\gamma}_{\omega})$ satisfying $\llbracket u^{\varepsilon} \rrbracket_{\alpha; \mathbf{R}^n} \leq C(n, \gamma) \llbracket u \rrbracket_{\alpha; \mathbf{R}^n}$ and $\lim_{\varepsilon \to 0} \llbracket u^{\varepsilon} - u \rrbracket_{\alpha/2; \mathbf{R}^n} = 0.$

Proof. By Hölder's inequality and Fubini's theorem we have

$$\mathbb{E}|[u^{\varepsilon}(x) - u^{\varepsilon}(x')|^{\gamma} = \mathbb{E}\left|\int_{B_{\varepsilon}} \varphi^{\varepsilon}(y)[u(x-y) - u(x'-y)] \,\mathrm{d}y\right|^{\gamma}$$

$$\leq |B_{\varepsilon}|^{\gamma-1} \int_{B_{\varepsilon}} |\varphi^{\varepsilon}(y)|^{\gamma} \,\mathrm{d}y\Big(\sup_{y \in B_{\varepsilon}} \mathbb{E}|u(x-y) - u(x'-y)|^{\gamma}\Big) \leq C[[u]]_{\alpha;\mathbf{R}^{n}}^{\gamma}|x-x'|^{\alpha\gamma},$$

thus $\llbracket u^{\varepsilon} \rrbracket_{\alpha;\mathbf{R}^n} \leq C \llbracket u \rrbracket_{\alpha;\mathbf{R}^n}$. Furthermore, by the definition of Hölder norms we can easily seen that

$$\llbracket u^{\varepsilon} - u \rrbracket_{\alpha/2;\mathbf{R}^n}^2 \leq \llbracket u^{\varepsilon} - u \rrbracket_{\alpha;\mathbf{R}^n} \llbracket u^{\varepsilon} - u \rrbracket_{0;\mathbf{R}^n} \leq C \llbracket u \rrbracket_{\alpha;\mathbf{R}^n} \llbracket u^{\varepsilon} - u \rrbracket_{0;\mathbf{R}^n},$$

so the proof is concluded by Lemma A.3.

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