## A FUNDAMENTAL MEAN-SQUARE CONVERGENCE THEOREM FOR SDES WITH LOCALLY LIPSCHITZ COEFFICIENTS AND ITS APPLICATIONS\*

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**Abstract.** A version of the fundamental mean-square convergence theorem is proved for stochastic differential equations (SDEs) in which coefficients are allowed to grow polynomially at infinity and which satisfy a one-sided Lipschitz condition. The theorem is illustrated on a number of particular numerical methods, including a special balanced scheme and fully implicit methods. The proposed special balanced scheme is explicit and its mean-square order of convergence is 1/2. Some numerical tests are presented.

**Key words.** SDEs with nonglobally Lipschitz coefficients, numerical integration of SDEs in the mean-square sense, balanced methods, fully implicit methods, strong convergence, almost sure convergence

AMS subject classifications. Primary, 60H35; Secondary, 65C30, 60H10

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**1. Introduction.** Let  $(\Omega, \mathcal{F}, P)$  be a complete probability space and  $\mathcal{F}_t^w$  be an increasing family of  $\sigma$ -subalgebras of  $\mathcal{F}$  induced by w(t) for  $0 \leq t \leq T$ , where  $(w(t), \mathcal{F}_t^w) = ((w_1(t), \dots, w_m(t))^\top, \mathcal{F}_t^w)$  is an *m*-dimensional standard Wiener process. We consider the system of Ito stochastic differential equations (SDEs)

(1.1) 
$$dX = a(t, X)dt + \sum_{r=1}^{m} \sigma_r(t, X)dw_r(t), \ t \in (t_0, T], \ X(t_0) = X_0,$$

where X, a,  $\sigma_r$  are d-dimensional column-vectors and  $X_0$  is independent of w. We suppose that any solution  $X_{t_0,X_0}(t)$  of (1.1) is regular on  $[t_0,T]$ . We recall [5] that a process is called regular if it is defined for all  $t_0 \leq t \leq T$ .

In traditional numerical analysis for SDEs [18, 12, 21] it is assumed that the SDEs coefficients are globally Lipschitz, which is a significant limitation taking into account that most of the models of applicable interest have coefficients which grow faster at infinity than a linear function. If the global Lipschitz condition is violated, the convergence of many usual numerical methods can disappear (see, e.g., [29, 6, 8, 22]). This has been the motivation for the recent interest in both theoretical support of existing numerical methods and developing new methods or approaches for solving SDEs under nonglobal Lipschitz assumptions on the coefficients.

In most SDEs applications (e.g., in molecular dynamics, financial engineering, and other problems of mathematical physics), one is interested in simulating averages  $\mathbb{E}\varphi(X(T))$  of the solution to SDEs—the task for which the weak-sense SDEs approximation is sufficient and effective [18, 21]. The problem with divergence of

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weak-sense schemes was addressed in [22] (see also [23]) for simulation of averages at finite time and also of ergodic limits when ensemble averaging is used. The concept of rejecting exploding trajectories proposed and justified in [22] allows us to use any numerical method for solving SDEs with nonglobally Lipschitz coefficients for estimating averages. Following this concept, we do not take into account the approximate trajectories X(t) which leave a sufficiently large ball  $S_R := \{x : |x| < R\}$  during the time T. See other approaches for resolving this problem in the context of computing averages, including the case of simulating ergodic limits via time averaging, e.g., in [29, 15, 1].

In this paper we deal with mean-square (strong) approximation of SDEs with nonglobal Lipschitz coefficients. Mean-square schemes have their own area of applicability (e.g., for simulating scenarios, visualization of stochastic dynamics, filtering; see further discussion on this in [12, 21, 9] and references therein). Furthermore, mean-square approximation is of theoretical interest and it also provides a guidance in constructing weak-sense schemes (see, e.g., [18, 12, 21]).

We note that in the case of weak approximation we often have to simulate large dimensional complicated stochastic systems using the Monte Carlo technique (or time averaging), which is typical for molecular dynamics applications, or we have to perform calculations on a daily basis, which is usual, e.g., in financial applications. Hence the cost per step of a weak numerical integrator should be low, which, in particular, essentially prohibits the use of implicit methods. In contrast, areas of applicability of mean-square schemes, as a rule, do not involve simulation of a large number of trajectories or over very long time periods and, consequently, there are more relaxed requirements on the cost per step of mean-square schemes and efficient and reliable implicit schemes have practical interest. There have been a number of recent works, including [8, 6, 10, 13, 11, 9, 14, 26, 7] (see also the references therein), where strong schemes for SDEs with nonglobal Lipschitz coefficients were considered. An extended literature review on this topic is available in [9].

In this paper we give a variant of the fundamental mean-square convergence theorem in the case of SDEs with nonglobal Lipschitz coefficients, which is analogous to Milstein's fundamental theorem for the global Lipschitz case [17] (see also [18, 21). More precisely, we assume that the SDEs coefficients can grow polynomially at infinity and satisfy a one-sided Lipschitz condition. The theorem is stated in section 2 and proved in Appendix A. Its corollary on almost sure convergence is also given. In section 2 we start discussion on applicability of the fundamental theorem, including its application to the drift-implicit Euler scheme, and thus establish its order of convergence. Strong convergence (but without order) of this scheme was proved for SDEs with nonglobal Lipschitz drift and diffusion in [13, 9] and more recently its convergence with order 1/2 was proved in [14]. A particular balanced method (see the class of balanced methods in [19, 21]) is proposed and its convergence with order 1/2 in the nonglobal Lipschitz setting is proved in section 3. Apparently, this is the first time mean-square convergence with an order has been proved for an explicit scheme under the conditions which allow polynomial growth of both drift and diffusion coefficients. In section 4 we revisit fully implicit (i.e., implicit both in drift and diffusion) mean-square schemes proposed and motivated by symplectic integration of stochastic Hamiltonian equations in [20] (see also [21]). In [20, 21] their convergence was proved for SDEs with globally Lipschitz coefficients. Here we relax these conditions as the drift is required to satisfy only a one-sided Lipschitz condition and be of not faster than polynomial growth at infinity. Some numerical experiments supporting our results are presented in section 5.

**2. Fundamental theorem.** Let  $X_{t_0,X_0}(t) = X(t), t_0 \le t \le T$ , be a solution of the system (1.1). We will assume the following.

Assumption 2.1. (i) The initial condition is such that

2.1) 
$$\mathbb{E}|X_0|^{2p} \le K < \infty \text{ for all } p \ge 1.$$

(ii) For a sufficiently large  $p_0 \ge 1$  there is a constant  $c_1 \ge 0$  such that for  $t \in [t_0, T]$ ,

$$(2.2) \ (x-y,a(t,x)-a(t,y)) + \frac{2p_0-1}{2} \sum_{r=1}^m |\sigma_r(t,x) - \sigma_r(t,y)|^2 \le c_1 |x-y|^2, \ x,y \in \mathbb{R}^d.$$

(iii) There exist  $c_2 \ge 0$  and  $\varkappa \ge 1$  such that for  $t \in [t_0, T]$ ,

(2.3) 
$$|a(t,x) - a(t,y)|^2 \le c_2(1+|x|^{2\varkappa-2}+|y|^{2\varkappa-2})|x-y|^2, x,y \in \mathbb{R}^d.$$

We note that (2.2) implies that

(2.4) 
$$(x, a(t, x)) + \frac{2p_0 - 3}{2} \sum_{r=1}^m |\sigma_r(t, x)|^2 \le c_0 + c_1' |x|^2, \quad t \in [t_0, T], \ x \in \mathbb{R}^d,$$

where  $c_0 = |a(t,0)|^2/2 + \frac{(2p_0-3)(2p_0-1)}{4} \sum_{r=1}^m |\sigma_r(t,0)|^2$  and  $c'_1 = c_1 + 1/2$ . The inequality (2.4) together with (2.1) is sufficient to ensure finiteness of moments [5]: there is K > 0,

(2.5) 
$$\mathbb{E}|X_{t_0,X_0}(t)|^{2p} < K(1+\mathbb{E}|X_0|^{2p}), \ 1 \le p \le p_0 - 1, \ t \in [t_0,T].$$

Also, (2.3) implies that

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(2.6) 
$$|a(t,x)|^2 \le c_3 + c'_2 |x|^{2\varkappa}, \quad t \in [t_0,T], \ x \in \mathbb{R}^d$$

where  $c_3 = 2|a(t,0))|^2 + 2c_2(\varkappa - 1)/\varkappa$  and  $c'_2 = 2c_2(1 + \varkappa)/\varkappa$ .

Introduce the one-step approximation  $X_{t,x}(t+h)$ ,  $t_0 \leq t < t+h \leq T$ , for the solution  $X_{t,x}(t+h)$  of (1.1), which depends on the initial point (t,x), a time step h, and  $\{w_1(\theta) - w_1(t), \ldots, w_m(\theta) - w_m(t), t \leq \theta \leq t+h\}$  and which is defined as follows:

(2.7) 
$$\bar{X}_{t,x}(t+h) = x + A(t,x,h;w_i(\theta) - w_i(t), \ i = 1,\dots,m, \ t \le \theta \le t+h).$$

Using the one-step approximation (2.7), we recurrently construct the approximation  $(X_k, \mathcal{F}_{t_k}), k = 0, \ldots, N, t_{k+1} - t_k = h_{k+1}, T_N = T$ :

(2.8) 
$$X_0 = X(t_0), \ X_{k+1} = \bar{X}_{t_k, \bar{X}_k}(t_{k+1}) \\ = X_k + A(t_k, X_k, h_{k+1}; w_i(\theta) - w_i(t_k), \ i = 1, \dots, m, \ t_k \le \theta \le t_{k+1})$$

The following theorem is a generalization of Milstein's fundamental theorem [17] (see also [18], [21, Chapter 1]) from the global to nonglobal Lipschitz case. It also has similarities with a strong convergence theorem in [6] proved for the case of nonglobal Lipschitz drift, global Lipschitz diffusion, and Euler-type schemes.

For simplicity, we will consider a uniform time step size, i.e.,  $h_k = h$  for all k. THEOREM 2.1. Suppose (i) Assumption 2.1 holds.

(ii) The one-step approximation  $\bar{X}_{t,x}(t+h)$  from (2.7) has the following orders of accuracy: for some  $p \ge 1$  there are  $\alpha \ge 1$ ,  $h_0 > 0$ , and K > 0 such that for arbitrary  $t_0 \le t \le T-h$ ,  $x \in \mathbb{R}^d$ , and all  $0 < h \le h_0$ ,

(2.9) 
$$|\mathbb{E}[X_{t,x}(t+h) - \bar{X}_{t,x}(t+h)]| \le K(1+|x|^{2\alpha})^{1/2}h^{q_1},$$

2.10) 
$$\left[\mathbb{E}|X_{t,x}(t+h) - \bar{X}_{t,x}(t+h)|^{2p}\right]^{1/(2p)} \le K(1+|x|^{2\alpha p})^{1/(2p)}h^{q_2}$$

with

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(2.11) 
$$q_2 \ge \frac{1}{2}, \ q_1 \ge q_2 + \frac{1}{2}.$$

(iii) The approximation  $X_k$  from (2.8) has finite moments, i.e., for some  $p \ge 1$  there are  $\beta \ge 1$ ,  $h_0 > 0$ , and K > 0 such that for all  $0 < h \le h_0$  and all k = 0, ..., N,

(2.12) 
$$\mathbb{E}|X_k|^{2p} < K(1 + \mathbb{E}|X_0|^{2p\beta})$$

Then for any N and k = 0, 1, ..., N the following inequality holds:

(2.13) 
$$\left[ \mathbb{E} |X_{t_0,X_0}(t_k) - \bar{X}_{t_0,X_0}(t_k)|^{2p} \right]^{1/(2p)} \le K(1 + \mathbb{E} |X_0|^{2\gamma p})^{1/(2p)} h^{q_2 - 1/2},$$

where K > 0 and  $\gamma \ge 1$  do not depend on h and k, i.e., the order of accuracy of the method (2.8) is  $q = q_2 - 1/2$ .

The theorem is proved in Appendix A and it uses the following lemma. LEMMA 2.2. Suppose Assumption 2.1 holds. For the representation

(2.14) 
$$X_{t,x}(t+\theta) - X_{t,y}(t+\theta) = x - y + Z_{t,x,y}(t+\theta),$$

we have for  $1 \le p \le (p_0 - 1)/\varkappa$ 

(2.15) 
$$\mathbb{E}|X_{t,x}(t+h) - X_{t,y}(t+h)|^{2p} \le |x-y|^{2p}(1+Kh),$$

(2.16) 
$$\mathbb{E} |Z_{t,x,y}(t+h)|^{2p} \le K(1+|x|^{2\varkappa-2}+|y|^{2\varkappa-2})^{p/2}|x-y|^{2p}h^{p}.$$

This lemma is proved in Appendix B. Theorem 2.1 has the following corollary.

COROLLARY 2.3. In the setting of Theorem 2.1 for  $p \ge 1/(2q)$  in (2.13), there is  $0 < \varepsilon < q$  and an a.s. finite random variable  $C(\omega) > 0$  such that

$$|X_{t_0,X_0}(t_k) - X_k| \le C(\omega) h^{q-\varepsilon},$$

i.e., the method (2.8) for (1.1) converges with order  $q - \varepsilon$  a.s.

The corollary is proved using the Borel–Cantelli-type of arguments (see, e.g., [3, 24]).

**2.1.** Discussion. In this section we make a number of observations concerning Theorem 2.1.

1. As a rule, it is not difficult to check the conditions (2.9)-(2.10) following the usual routine calculations as in the global Lipschitz case [18, 12, 21]. We note that in order to achieve the optimal  $q_1$  and  $q_2$  in (2.9)-(2.10) additional assumptions on smoothness of a(t, x) and  $\sigma_r(t, x)$  are usually needed.

In contrast to the conditions (2.9)-(2.10), checking the condition (2.12) on moments of a method  $X_k$  is often rather difficult. In the case of global Lipschitz coefficients, boundedness of moments of  $X_k$  is just direct implication of the boundedness of moments of the SDEs solution and the one-step properties of the method (see [21, Lemma 1.1.5]). There is no result of this type in the case of the nonglobal Lipschitz SDEs and each scheme requires a special consideration. For a number of strong schemes boundedness of moments in nonglobal Lipschitz cases were proved (see, e.g., [8, 6, 10, 9, 29]). In section 3 we show boundedness of moments for a balanced method and in section 4 for fully implicit methods.

Roughly speaking, Theorem 2.1 says that if moments of  $X_k$  are bounded and the scheme was proved to be convergent with order q in the global Lipschitz case, then the scheme has the same convergence order q in the considered nonglobal Lipschitz case.

2. Assumptions and the statement of Theorem 2.1 include the famous fundamental theorem of Milstein [17] proved under the global conditions on the SDEs coefficients. (Of course, as discussed in the previous point, this case does not need the assumption (2.12).)

3. Consider the drift-implicit scheme [21, p. 30]

(2.17) 
$$X_{k+1} = X_k + a(t_{k+1}, X_{k+1})h + \sum_{r=1}^m \sigma_r(t_k, X_k)\xi_{rk}\sqrt{h},$$

where  $\xi_{rk} = (w_r(t_{k+1}) - w_r(t_k))/\sqrt{h}$  are Gaussian  $\mathcal{N}(0, 1)$  independent and identically distributed (i.i.d.) random variables. Assume that the coefficients a(t, x) and  $\sigma_r(t, x)$ have continuous first-order partial derivatives in t and the coefficient a(t, x) also has continuous first-order partial derivatives in  $x^i$  and that all these derivatives and the coefficients themselves satisfy inequalities of the form (2.3). It is not difficult to show that the one-step approximation corresponding to (2.17) satisfies (2.9) and (2.10)with  $q_1 = 2$  and  $q_2 = 1$ , respectively. Its boundedness of moments, in particular, under the condition (2.4) for time steps  $h \leq 1/(2c_1)$ , is proved in [9]. Then, due to Theorem 2.1, (2.17) converges with mean-square order q = 1/2. (Note that for q = 1/2, it is sufficient to have  $q_1 = 3/2$ , which can be obtained under lesser smoothness of a.) Further, in the case of additive noise (i.e.,  $\sigma_r(t, x) = \sigma_r(t), r = 1, \ldots, m$ ),  $q_1 = 2$  and  $q_2 = 3/2$  and (2.17) converges with mean-square order 1 due to Theorem 2.1. We note that convergence of (2.17) with order 1/2 in the global Lipschitz case is well known [18, 12, 21]; in the case of nonglobal Lipschitz drift and global Lipschitz diffusion it was proved in [8, 6] (see also related results in [3, 29]); and under Assumption 2.1 strong convergence of (2.17) without order was proved in [13, 9] and more recently its strong order 1/2 was independently established in [14].

4. Due to the bound (2.5) on the moments of the solution X(t), it would be natural to require that  $\beta$  in (2.12) be equal to 1. Indeed, (2.12) with  $\beta = 1$  holds for the drift-implicit method (2.17) [9] and for fully implicit methods (see section 4). However, this is not the case for tamed-type methods (see [10]) or the balanced method from section 3.

5. The constant K in (2.13) depends on  $p, t_0, T$  as well as on the SDEs coefficients. The constant  $\gamma$  in (2.13) depends on  $\alpha$ ,  $\beta$ , and  $\varkappa$ .

6. Let us illustrate Assumption 2.1(ii) on a one-dimensional SDE:  $dX = -\mu X |X|^{r_1-1} dt + \lambda X^{r_2} dw$  with  $\mu$ ,  $\lambda > 0$ ,  $r_1 \ge 1$ , and  $r_2 \ge 1$ . If  $r_1 + 1 > 2r_2$  or  $r_1 = r_2 = 1$ , then (2.2) is valid for any  $p_0 \ge 1$ . If  $r_1 + 1 = 2r_2$  and  $r_1 > 1$ , then (2.2) is valid for  $1 \le p_0 \le \mu/\lambda^2 + 1/2$ .

**3.** A balanced method. In this section we propose a particular balanced scheme from the class of balanced methods introduced in [19] (see also [21]) and prove its mean-square convergence with order 1/2 using Theorem 2.1. As far as we know, this variant of balanced schemes has not been considered before. In section 5 we test the balanced scheme on a model problem and demonstrate that it is more efficient than the tamed scheme (5.2) (see section 5) from [9]. We also note that it was mentioned in [9] that a balanced scheme suitable for the nonglobal Lipschitz case could potentially be derived.

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Consider the following balanced-type scheme for (1.1):

(3.1) 
$$X_{k+1} = X_k + \frac{a(t_k, X_k)h + \sum_{r=1}^m \sigma_r(t_k, X_k)\xi_{rk}\sqrt{h}}{1 + h|a(t_k, X_k)| + \sqrt{h}\sum_{r=1}^m |\sigma_r(t_k, X_k)\xi_{rk}|}$$

where  $\xi_{rk}$  are Gaussian  $\mathcal{N}(0, 1)$  i.i.d. random variables.

We prove two lemmas which show that the scheme (3.1) satisfies the conditions of Theorem 2.1. The first lemma is on boundedness of moments, which uses a stopping time technique (see also, e.g., [22, 9]).

LEMMA 3.1. Suppose Assumption 2.1 holds with sufficiently large  $p_0$ . For all natural N and all k = 0, ..., N the following inequality holds for moments of the scheme (3.1):

(3.2) 
$$\mathbb{E}|X_k|^{2p} \le K(1 + \mathbb{E}|X_0|^{2p\beta}), \quad 1 \le p \le \frac{p_0 - 1}{4(3\varkappa - 2)} - \frac{1}{2}$$

with some constants  $\beta \geq 1$  and K > 0 independent of h and k.

*Proof.* In the proof we shall use the letter K to denote various constants which are independent of h and k. We note in passing that the case  $\varkappa = 1$  (i.e., when a(t, x) is globally Lipschitz) is trivial.

The following elementary consequence of the inequalities (2.4) and (2.6) will be used in the proof: there exits a constant K > 0 such that

(3.3) 
$$\sum_{r=1}^{m} |\sigma_r(t,x)|^2 \le K(1+|x|^{2\varkappa}).$$

We observe from (3.1) that

(3.4) 
$$|X_{k+1}| \le |X_k| + 1 \le |X_0| + (k+1)$$

Let R > 0 be a sufficiently large number. Introduce the events

(3.5) 
$$\tilde{\Omega}_{R,k} := \{\omega : |X_l| \le R, \ l = 0, \dots, k\}$$

and their compliments  $\tilde{\Lambda}_{R,k}$ . We first prove the lemma for integer  $p \geq 1$ . We have

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$$\begin{split} & \mathbb{E}\chi_{\tilde{\Omega}_{R,k+1}}(\omega)|X_{k+1}|^{2p} \leq \mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|X_{k+1}|^{2p} \\ & = \mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|(X_{k+1} - X_k) + X_k|^{2p} \leq \mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|X_k|^{2p} + \mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|X_k|^{2p-2} \\ & \times \left[2p(X_k, X_{k+1} - X_k) + p(2p-1)|X_{k+1} - X_k|^2\right] \\ & + K\sum_{l=3}^{2p} \mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|X_k|^{2p-l}|X_{k+1} - X_k|^l. \end{split}$$

Consider the second term in the right-hand side of (3.6):

$$\begin{split} & \mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega) \left|X_{k}\right|^{2p-2} \left[2p(X_{k}, X_{k+1} - X_{k}) + p(2p-1)|X_{k+1} - X_{k}|^{2}\right] \\ &= 2p\mathbb{E}\left(\chi_{\tilde{\Omega}_{R,k}}(\omega) \left|X_{k}\right|^{2p-2} \mathbb{E}\left[\left(X_{k}, \frac{a(t_{k}, X_{k})h + \sum_{r=1}^{m} \sigma_{r}(t_{k}, X_{k})\xi_{rk}\sqrt{h}}{1 + h|a(t_{k}, X_{k})| + \sqrt{h}\sum_{r=1}^{m} |\sigma_{r}(t_{k}, X_{k})\xi_{rk}|}\right) \\ &+ \frac{2p-1}{2}\left|\frac{a(t_{k}, X_{k})h + \sum_{r=1}^{m} \sigma_{r}(t_{k}, X_{k})\xi_{rk}\sqrt{h}}{1 + h|a(t_{k}, X_{k})| + \sqrt{h}\sum_{r=1}^{m} |\sigma_{r}(t_{k}, X_{k})\xi_{rk}|}\right|^{2} \left|\mathcal{F}_{t_{k}}\right|\right). \end{split}$$

Since  $\mathbb{E}\xi_{rk}\prod_{j=1}^{m} |\xi_{jk}|^{\alpha_j} = 0$  for all r and any  $\alpha_j \geq 0, j = 1, \ldots m$ , and  $\xi_{rk}$  are independent of  $\mathcal{F}_{t_k}$ , we obtain

(3.8) 
$$\chi_{\tilde{\Omega}_{R,k}} \mathbb{E}\left[\frac{\sum_{r=1}^{m} \sigma_r(t_k, X_k)\xi_{rk}\sqrt{h}}{1+h|a(t_k, X_k)| + \sqrt{h}\sum_{r=1}^{m} |\sigma_r(t_k, X_k)\xi_{rk}|} \middle| \mathcal{F}_{t_k}\right]$$
$$= \chi_{\tilde{\Omega}_{R,k}} \sum_{r=1}^{m} \mathbb{E}\left[\sigma_r(t_k, X_k)\xi_{rk}\sqrt{h}\sum_{i=0}^{\infty} (-1)^i \left[h|a(t_k, X_k)| + \sqrt{h}\sum_{r=1}^{m} |\sigma_r(t_k, X_k)\xi_{rk}|\right]^i \middle| \mathcal{F}_{t_k}\right] = 0.$$

Using that  $\mathbb{E}\xi_{rk}\xi_{lk}\prod_{j=1}^{m} |\xi_{jk}|^{\alpha_j} = 0$  for  $l \neq r$  and any  $\alpha_j \geq 0, j = 1, \ldots, m$ , we analogously get for  $l \neq r$ 

(3.9) 
$$\chi_{\tilde{\Omega}_{R,k}} \mathbb{E}\left[\frac{\sigma_r(t_k, X_k)\xi_{rk}\sqrt{h}\sigma_l(t_k, X_k)\xi_{lk}\sqrt{h}}{(1+h|a(t_k, X_k)| + \sqrt{h}\sum_{r=1}^m |\sigma_r(t_k, X_k)\xi_{rk}|)^2}\right|\mathcal{F}_{t_k}\right] = 0.$$

Then the conditional expectation in (3.7) becomes

$$(3.10) A := \chi_{\bar{\Omega}_{R,k}} \mathbb{E} \left[ \left( X_k, \frac{a(t_k, X_k)h + \sum_{r=1}^m \sigma_r(t_k, X_k)\xi_{rk}\sqrt{h}}{1 + h|a(t_k, X_k)| + \sqrt{h}\sum_{r=1}^m |\sigma_r(t_k, X_k)\xi_{rk}\sqrt{h}} \right|^2 |\mathcal{F}_{t_k} \right] \\ + \frac{2p-1}{2} \left| \frac{a(t_k, X_k)h + \sum_{r=1}^m \sigma_r(t_k, X_k)\xi_{rk}\sqrt{h}}{1 + h|a(t_k, X_k)| + \sqrt{h}\sum_{r=1}^m |\sigma_r(t_k, X_k)\xi_{rk}|} \right|^2 |\mathcal{F}_{t_k} \right] \\ = \chi_{\bar{\Omega}_{R,k}} \mathbb{E} \left[ \frac{(X_k, a(t_k, X_k)h)}{1 + h|a(t_k, X_k)| + \sqrt{h}|\sum_{r=1}^m \sigma_r(t_k, X_k)\xi_{rk}|} \right. \\ + \frac{2p-1}{2} \frac{a^2(t_k, X_k)h^2 + h\sum_{r=1}^m (\sigma_r(t_k, X_k)\xi_{rk})^2}{(1 + h|a(t_k, X_k)| + \sqrt{h}|\sum_{r=1}^m \sigma_r(t_k, X_k)\xi_{rk}|}} \right] \\ \leq \chi_{\bar{\Omega}_{R,k}} \mathbb{E} \left[ \frac{(X_k, a(t_k, X_k)h)}{1 + h|a(t_k, X_k)| + \sqrt{h}\sum_{r=1}^m |\sigma_r(t_k, X_k)\xi_{rk}|} \right. \\ + \frac{2p-1}{2} \frac{h\sum_{r=1}^m |\sigma_r(t_k, X_k)|^2 \xi_{rk}^2}{1 + h|a(t_k, X_k)| + \sqrt{h}\sum_{r=1}^m |\sigma_r(t_k, X_k)\xi_{rk}|}} \right| \mathcal{F}_{t_k} \right] \\ + \frac{2p-1}{2} \chi_{\bar{\Omega}_{R,k}} a^2(t_k, X_k)h^2 \\ = \chi_{\bar{\Omega}_{R,k}} \mathbb{E} \left[ \frac{(X_k, a(t_k, X_k)h) + \frac{2p-1}{2}h\sum_{r=1}^m |\sigma_r(t_k, X_k)\xi_{rk}|}{1 + h|a(t_k, X_k)| + \sqrt{h}\sum_{r=1}^m |\sigma_r(t_k, X_k)\xi_{rk}|}} \right| \mathcal{F}_{t_k} \right] \\ + \frac{2p-1}{2} \chi_{\bar{\Omega}_{R,k}} a^2(t_k, X_k)h^2.$$

Using (2.4) and (2.6), we obtain

(3.11) 
$$A \leq c_0 h + c_1' |X_k|^2 h \chi_{\tilde{\Omega}_{R,k}} + \frac{2p-1}{2} h \chi_{\tilde{\Omega}_{R,k}} \sum_{r=1}^m |\sigma_r(t_k, X_k)|^2 \\ \times \mathbb{E}\left[\frac{(\xi_{rk}^2 - 1)}{1 + h|a(t_k, X_k)| + \sqrt{h} \sum_{r=1}^m |\sigma_r(t_k, X_k)\xi_{rk}|} \middle| \mathcal{F}_{t_k} \right] \\ + K h^2 + K \chi_{\tilde{\Omega}_{R,k}} |X_k|^{2\varkappa} h^2.$$

Since  $\mathbb{E}(\xi_{rk}^2 - 1) = 0$ , moments of  $\xi_{rk}$  are bounded, and  $\xi_{rk}$  are independent of  $\mathcal{F}_{t_k}$ , we obtain for the expectation in the second term in (3.11)

$$(3.12) \quad \chi_{\tilde{\Omega}_{R,k}} \mathbb{E} \left[ \frac{(\xi_{rk}^{2} - 1)}{1 + h|a(t_{k}, X_{k})| + \sqrt{h} \sum_{r=1}^{m} |\sigma_{r}(t_{k}, X_{k})\xi_{rk}|} \middle| \mathcal{F}_{t_{k}} \right] \\ = \chi_{\tilde{\Omega}_{R,k}} \mathbb{E} \left[ \frac{(\xi_{rk}^{2} - 1)}{1 + h|a(t_{k}, X_{k})| + \sqrt{h} \sum_{r=1}^{m} |\sigma_{r}(t_{k}, X_{k})\xi_{rk}|} - (\xi_{rk}^{2} - 1) \middle| \mathcal{F}_{t_{k}} \right] \\ = -\chi_{\tilde{\Omega}_{R,k}} \mathbb{E} \left[ (\xi_{rk}^{2} - 1) \frac{h|a(t_{k}, X_{k})| + \sqrt{h} \sum_{r=1}^{m} |\sigma_{l}(t_{k}, X_{k})\xi_{lk}|}{1 + h|a(t_{k}, X_{k})| + \sqrt{h} \sum_{l=1}^{m} |\sigma_{r}(t_{k}, X_{k})\xi_{lk}|} \middle| \mathcal{F}_{t_{k}} \right] \\ \leq \chi_{\tilde{\Omega}_{R,k}} \mathbb{E} \left[ |\xi_{rk}^{2} - 1| \left( h|a(t_{k}, X_{k})| + \sqrt{h} \sum_{r=1}^{m} |\sigma_{l}(t_{k}, X_{k})||\xi_{lk}| \right) \middle| \mathcal{F}_{t_{k}} \right] \\ \leq \chi_{\tilde{\Omega}_{R,k}} K \left( h|a(t_{k}, X_{k})| + \sqrt{h} \sum_{r=1}^{m} |\sigma_{r}(t_{k}, X_{k})| \right).$$

Using (2.6) and (3.3), we get from (3.11)-(3.12)

$$(3.13) A \leq c_0 h + c_1' \chi_{\tilde{\Omega}_{R,k}} |X_k|^2 h + Kh \chi_{\tilde{\Omega}_{R,k}} \sum_{r=1}^m |\sigma_r(t_k, X_k)|^2 \\ \times \left[ h |a(t_k, X_k)| + \sqrt{h} \sum_{r=1}^m |\sigma_r(t_k, X_k)| \right] + Kh^2 + K \chi_{\tilde{\Omega}_{R,k}} |X_k|^{2\varkappa} h^2 \\ \leq \chi_{\tilde{\Omega}_{R,k}} Kh (1 + |X_k|^2 + |X_k|^{2\varkappa} h + |X_k|^{3\varkappa} h^{1/2}) \\ \leq \chi_{\tilde{\Omega}_{R,k}} Kh (1 + |X_k|^2 + |X_k|^{3\varkappa} h^{1/2}).$$

Now consider the last term in (3.6):

(3.14) 
$$\mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega) |X_k|^{2p-l} |X_{k+1} - X_k|^l \\ \leq K \mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega) |X_k|^{2p-l} \left[ h^l |a(t_k, X_k)|^l + h^{l/2} \sum_{r=1}^m |\sigma_r(t_k, X_k)|^l |\xi_{rk}|^l \right] \\ \leq K \mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega) |X_k|^{2p-l} h^{l/2} \left[ 1 + |X_k|^{l\varkappa} \right],$$

where we used (2.6) and (3.3) again as well as the fact that  $\chi_{\tilde{\Omega}_{R,k}}(\omega)$  and  $X_k$  are  $\mathcal{F}_{t_k}$ -measurable while  $\xi_{rk}$  are independent of  $\mathcal{F}_{t_k}$ .

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Combining (3.6), (3.7), (3.10), (3.13), and (3.14), we obtain

(3.15)

$$\begin{split} & \mathbb{E}\chi_{\tilde{\Omega}_{R,k+1}}(\omega)|X_{k+1}|^{2p} \\ & \leq \mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|X_{k}|^{2p} + Kh\mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|X_{k}|^{2p-2}\left[1 + |X_{k}|^{2} + |X_{k}|^{3\varkappa}h^{1/2}\right] \\ & + K\sum_{l=3}^{2p}\mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|X_{k}|^{2p-l}h^{l/2}\left[1 + |X_{k}|^{l\varkappa}\right] \\ & \leq \mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|X_{k}|^{2p} + Kh\mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|X_{k}|^{2p} + K\sum_{l=2}^{2p}\mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|X_{k}|^{2p-l}h^{l/2} \\ & + Kh^{3/2}\mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|X_{k}|^{2p-2+3\varkappa} + Kh\sum_{l=3}^{2p}\mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega)|X_{k}|^{2p+l(\varkappa-1)}h^{l/2-1}. \end{split}$$

Choosing

(3.16) 
$$R = R(h) = h^{-1/(6\varkappa - 4)},$$

we get, for  $l = 3, \ldots, 2p$ ,  $\mathbb{E}\chi_{\tilde{\Omega}_{R,k}}(\omega) |X_k|^{2p-2+3\varkappa} h^{l/2-1} \leq \chi_{\tilde{\Omega}_{R(h),k}}(\omega) |X_k|^{2p}$  and  $\chi_{\tilde{\Omega}_{R(h),k}}(\omega) |X_k|^{2p+l(\varkappa-1)} h^{l/2-1} \leq \chi_{\tilde{\Omega}_{R(h),k}}(\omega) |X_k|^{2p}$ , and hence we rewrite (3.15) as

$$\begin{aligned} & \mathbb{E}\chi_{\tilde{\Omega}_{R(h),k+1}}(\omega)|X_{k+1}|^{2p} \\ & \leq \mathbb{E}\chi_{\tilde{\Omega}_{R(h),k}}(\omega)|X_{k}|^{2p} + Kh\mathbb{E}\chi_{\tilde{\Omega}_{R(h),k}}(\omega)|X_{k}|^{2p} + K\sum_{l=1}^{p}\mathbb{E}\chi_{\tilde{\Omega}_{R(h),k}}(\omega)|X_{k}|^{2(p-l)}h^{l} \\ & \leq \mathbb{E}\chi_{\tilde{\Omega}_{R(h),k}}(\omega)|X_{k}|^{2p} + Kh\mathbb{E}\chi_{\tilde{\Omega}_{R(h),k}}(\omega)|X_{k}|^{2p} + Kh, \end{aligned}$$

where in the last line we have used Young's inequality. From here, we get by Gronwall's inequality that

(3.18) 
$$\mathbb{E}\chi_{\tilde{\Omega}_{R(h),k}}(\omega)|X_k|^{2p} \le K(1+\mathbb{E}|X_0|^{2p}),$$

where R(h) is from (3.16) and K does not depend on k and h but it depends on p.

It remains to estimate  $\mathbb{E}\chi_{\tilde{\Lambda}_{R(h),k}}(\omega)|X_k|^{2p}.$  We have

$$\chi_{\tilde{\Lambda}_{R,k}} = 1 - \chi_{\tilde{\Omega}_{R,k}} = 1 - \chi_{\tilde{\Omega}_{R,k-1}} \chi_{|X_k| \le R} = \chi_{\tilde{\Lambda}_{R,k-1}} + \chi_{\tilde{\Omega}_{R,k-1}} \chi_{|X_k| > R}$$
$$= \dots = \sum_{l=0}^k \chi_{\tilde{\Omega}_{R,l-1}} \chi_{|X_l| > R},$$

where we put  $\chi_{\tilde{\Omega}_{R,-1}} = 1$ . Then, using (3.4), (3.18), (2.1), and the Cauchy–Bunyakovsky–Schwarz and Markov inequalities, we obtain

(3.19)

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$$\begin{split} \mathbb{E}\chi_{\tilde{\Lambda}_{R(h),k}}(\omega)|X_{k}|^{2p} &= \mathbb{E}\sum_{l=0}^{k}|X_{k}|^{2p}\chi_{\tilde{\Omega}_{R(h),l-1}}\chi_{|X_{l}|>R(h)} \\ &\leq \left(\mathbb{E}|X_{0}+k|^{4p}\right)^{1/2}\sum_{l=0}^{k}\left(\mathbb{E}\left[\chi_{\tilde{\Omega}_{R(h),l-1}|X_{l}|>R(h)}\right]\right)^{1/2} \\ &= \left(\mathbb{E}|X_{0}+k|^{4p}\right)^{1/2}\sum_{l=0}^{k}\left(P(\chi_{\tilde{\Omega}_{R(h),l-1}}|X_{l}|>R)\right)^{1/2} \\ &\leq \left(\mathbb{E}|X_{0}+k|^{4p}\right)^{1/2}\sum_{l=0}^{k}\frac{\left(\mathbb{E}(\chi_{\tilde{\Omega}_{R(h),l-1}}|X_{l}|^{2(2p+1)(6\varkappa-4)})\right)^{1/2}}{R(h)^{(2p+1)(6\varkappa-4)}} \\ &\leq K\left(\mathbb{E}|X_{0}+k|^{4p}\right)^{1/2}\left(\mathbb{E}(1+|X_{0}|^{2(2p+1)(6\varkappa-4)})\right)^{1/2}kh^{2p+1} \\ &\leq K(1+\mathbb{E}|X_{0}|^{4p+2(2p+1)(6\varkappa-4)})^{1/2}, \end{split}$$

which together with (3.18) implies (3.2) for integer  $p \ge 1$ . Then, by Jensen's inequality, (3.2) holds for noninteger p as well.

The next lemma gives estimates for the one-step error of the balanced scheme (3.1).

LEMMA 3.2. Assume that (2.5) holds. Assume that the coefficients a(t, x) and  $\sigma_r(t, x)$  have continuous first-order partial derivatives in t and that these derivatives and the coefficients satisfy inequalities of the form (2.3). Then the scheme (3.1) satisfies the inequalities (2.9) and (2.10) with  $q_1 = 3/2$  and  $q_2 = 1$ , respectively.

The proof of this lemma is given in Appendix C. Lemmas 3.1 and 3.2 and Theorem 2.1 imply the following result.

PROPOSITION 3.3. Under the assumptions of Lemmas 3.1 and 3.2 the balanced scheme (3.1) has mean-square order 1/2, i.e., for it the inequality (2.13) holds with  $q = q_2 - 1/2 = 1/2$ .

*Remark* 3.1. In the additive noise case the mean-square order of the balanced scheme (3.1) does not improve  $(q_1 \text{ and } q_2 \text{ remain } 3/2 \text{ and } 1$ , respectively).

4. Fully implicit schemes. Fully implicit (i.e., implicit in both drift and diffusion coefficients) mean-square schemes were proposed in [20] (see also [21, Chapter 1]), where their convergence was proved under global Lipschitz conditions. We recall that in [20, 21] introduction of fully implicit methods was motivated by consideration of Hamiltonian equations with multiplicative noise since in the case of a general Hamiltonian system only implicit schemes can be symplectic. (We also note that in the case of deterministic general Hamiltonian systems symplectic Runge–Kutta methods are all implicit [28].) They have other interesting features too. The most remarkable scheme of fully implicit methods is the midpoint scheme, which complies with the Stratonovich calculus without any need for differentiating the diffusion coefficient and which is of the first mean-square order for Stratonovich SDEs with commutative noise (see [20, 21] and also Remark 4.2 below). The midpoint scheme is known for its good geometric integration features not only in the context of Hamiltonian systems but also, e.g., in the case of stochastic Landau–Lifshitz equations (see [16] and references therein).

Here we analyze fully implicit schemes under the following assumptions, which are stronger with respect to the diffusion coefficient than Assumption 2.1 used in sections 2 and 3.

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Assumption 4.1. (i) The initial condition is such that

(4.1) 
$$\mathbb{E}|X_0|^{2p} \le K < \infty \text{ for all } p \ge 1$$

(ii) There exists a constant  $c_1 \ge 0$  such that

(4.2) 
$$(x - y, a(t, x) - a(t, y)) \le c_1 |x - y|^2, t \in [t_0, T], x, y \in \mathbb{R}^d.$$

(iii) There exist  $c_2 \ge 0$  and  $\varkappa \ge 1$  such that

(4.3) 
$$|a(t,x) - a(t,y)|^2 \le c_2(1+|x|^{2\varkappa-2}+|y|^{2\varkappa-2})|x-y|^2, t \in [t_0,T], x,y \in \mathbb{R}^d.$$

(iv) The coefficients  $\sigma_r(t, x)$  have continuous bounded first-order spatial derivatives so that for  $t \in [t_0, T]$ , there are constants  $L_1 \ge 0$  and  $L_2 \ge 0$ :

(4.4) 
$$|\nabla \sigma_r(t,x)| \le L_1, \ r = 1, \dots, m, \ x \in \mathbb{R}^d,$$

and

$$(4.5) \quad |\nabla \sigma_r(t,x)\sigma_r(t,x) - \nabla \sigma_r(t,y)\sigma_r(t,y)| \le L_2|x-y|, \ r=1,\ldots,m, \ x,y \in \mathbb{R}^d.$$

In proofs which follow we will need some implications of Assumption 4.1. The condition (4.2) implies that there is  $c \ge 0$ 

(4.6) 
$$(x, a(t, x)) \le c(1 + |x|^2), \quad t \in [t_0, T], \ x \in \mathbb{R}^d.$$

It follows from (4.4) that

(4.7) 
$$|\sigma_r(t,x) - \sigma_r(t,y)| \le L_1 |x-y|, \ t \in [t_0,T], \ x,y \in \mathbb{R}^d,$$

and hence

(4.8) 
$$|\sigma_r(t,x)| \le L_1|x| + L_0,$$

where  $L_0 = \max_{t \in [t_0,T]} |\sigma_r(t,0)|$ . Further, there is  $L \ge 0$ :

(4.9) 
$$|\nabla \sigma_r(t,x)\sigma_r(t,x)| \le L(1+|x|), \quad t \in [t_0,T], \ x \in \mathbb{R}^d,$$

and

(4.10) 
$$|\sigma_r(t,x)|^2 \le L(1+|x|^2), \quad t \in [t_0,T], \ x \in \mathbb{R}^d.$$

For definiteness, we consider the following one-parametric family of methods for (1.1) from the broader class of fully implicit schemes of [20, 21]:

$$\begin{aligned} X_{k+1} &= X_k + a(t_{k+\lambda}, (1-\lambda)X_k + \lambda X_{k+1})h \\ &- \lambda \sum_{r=1}^m \sum_{j=1}^d \frac{\partial \sigma_r}{\partial x^j} (t_{k+\lambda}, (1-\lambda)X_k + \lambda X_{k+1}) \sigma_r^j (t_{k+\lambda}, (1-\lambda)X_k + \lambda X_{k+1})h \\ &+ \sum_{r=1}^m \sigma_r (t_{k+\lambda}, (1-\lambda)X_k + \lambda X_{k+1}) (\zeta_{rh})_k \sqrt{h}, \end{aligned}$$

where  $0 \leq \lambda \leq 1$ ,  $t_{k+\lambda} = t_k + \lambda h$ , and  $(\zeta_{rh})_k$  are i.i.d. random variables so that

(4.12) 
$$\zeta_{h} = \begin{cases} \xi, \ |\xi| \le A_{h}, \\ A_{h}, \ \xi > A_{h}, \\ -A_{h}, \ \xi < -A_{h} \end{cases}$$

with  $\xi \sim \mathcal{N}(0,1)$  and  $A_h = \sqrt{2l |\ln h|}$  with  $l \ge 1$ . We recall [21, Lemma 1.3.4] that

(4.13) 
$$E(\xi^2 - \zeta_h^2) = (1 + 2\sqrt{2l|\ln h|})h^l.$$

Remark 4.1. Three choices of  $\lambda$  are most notable:  $\lambda = 0$  gives the explicit Euler scheme which is divergent [6, 11] in the considered setting;  $\lambda = 1$  gives the fully implicit Euler scheme; and  $\lambda = 1/2$  corresponds to the midpoint rule, which in application to a system of Stratonovich SDEs is derivative free [21, p. 45].

Now we will study properties of the method (4.11). Consider the one-step approximations corresponding to (4.11),

(4.14) 
$$\bar{X} = \bar{X}^{\lambda} = x + a(t + \lambda h, U^{\lambda})h - \lambda \sum_{r=1}^{m} \sum_{j=1}^{d} \frac{\partial \sigma_{r}}{\partial x^{j}} (t + \lambda h, U^{\lambda}) \sigma_{r}^{j} (t + \lambda h, U^{\lambda})h + \sum_{r=1}^{m} \sigma_{r} (t + \lambda h, U^{\lambda}) \zeta_{rh} \sqrt{h},$$

where

(4.15) 
$$U = U^{\lambda} := (1 - \lambda)x + \lambda \bar{X}^{\lambda}.$$

Note that

$$(4.16) U^{\lambda} = x + \lambda a(t + \lambda h, U^{\lambda}) - \lambda^{2} \sum_{r=1}^{m} \sum_{j=1}^{d} \frac{\partial \sigma_{r}}{\partial x^{j}} (t + \lambda h, U^{\lambda}) \sigma_{r}^{j} (t + \lambda h, U^{\lambda}) h$$
$$+ \lambda \sum_{r=1}^{m} \sigma_{r} (t + \lambda h, U^{\lambda}) \zeta_{rh} \sqrt{h}.$$

LEMMA 4.1. Let  $0 < \lambda \leq 1$ . Assume that Assumption 4.1 holds. For an arbitrary  $0 < \varepsilon < 1$ , find  $h_0 > 0$  such that

(4.17) 
$$\lambda \left[ h_0 c_1 + m \lambda L_2 h_0 + m L_1 \sqrt{2lh_0 |\ln h_0|} \right] = 1 - \varepsilon.$$

Then (4.14) for any  $0 < h \le h_0$  has the unique solution  $\overline{X}$  which satisfies the inequalities for some K > 0:

(4.18) 
$$|\bar{X} - x| \le K(1 + |x|^{\varkappa})h + K(1 + |x|)\sqrt{h|\ln h|}$$

and

$$(4.19) \quad |\bar{X}|^2 \le \frac{16}{3\varepsilon^2\lambda} (L_0+1)\sqrt{2lh|\ln h|} + \frac{4}{\lambda^2} \left[ (1-\lambda)^2 + \frac{4}{3\varepsilon^2} \right] |x|^2, \ t \in [t_0,T], \ x \in \mathbb{R}^d.$$

Proof. Let

(4.20) 
$$\tilde{a}(t,x) = a(t,x) - \lambda \sum_{r=1}^{m} \sum_{j=1}^{d} \frac{\partial \sigma_r}{\partial x^j}(t,x) \sigma_r^j(t,x).$$

For any fixed  $\lambda$ , t,  $\zeta_{rh}$ , and h, we introduce the function

$$\psi(z) = z - \lambda \tilde{a}(t + \lambda h, z)h - \lambda \sum_{r=1}^{m} \sigma_r(t + \lambda h, z)\zeta_{rh}\sqrt{h},$$

which is continuous in z due to our assumptions. Equation (4.14) can be written as

(4.21) 
$$\psi(U^{\lambda}) = x.$$

Using (4.2), (4.5), and (4.7), we obtain

4.22) 
$$(z - y, \psi(z) - \psi(y)) \\ \ge |z - y|^2 - h\lambda c_1 |z - y|^2 - hm\lambda^2 L_2 |z - y|^2 - m\lambda L_1 \sqrt{2lh|\ln h|} |z - y|^2 \\ = \left(1 - \lambda \left[hc_1 + m\lambda L_2 h + mL_1 \sqrt{2lh|\ln h|}\right]\right) |z - y|^2 \ge \varepsilon |z - y|^2 > 0,$$

i.e.,  $\psi(z)$  is a uniformly monotone function for  $h \leq h_0$ . This implies (see, e.g., [27, Theorem 6.4.4, p. 167]) that (4.14) has a unique solution.

We obtain from (4.21) and (4.22)

$$\varepsilon |U|^{2} \leq (U, \psi(U) - \psi(0)) = (U, x - \psi(0))$$
  
$$\leq \frac{\varepsilon}{4} |U|^{2} + \frac{2}{\varepsilon} |x|^{2} + \frac{2}{\varepsilon} |\psi(0)|^{2} \leq \frac{\varepsilon}{4} |U|^{2} + \frac{2}{\varepsilon} |x|^{2} + \frac{2\lambda(L_{0} + 1)\sqrt{2lh|\ln h|}}{\varepsilon}$$

and hence

(4.23) 
$$|U|^2 \le \frac{8}{3\varepsilon^2} (\lambda (L_0 + 1)\sqrt{2lh|\ln h|} + |x|^2),$$

which implies (4.19).

Further, it follows from (4.15), (4.21), and (4.22) that

$$\begin{split} \lambda \varepsilon |\bar{X} - x|^2 &= \varepsilon |U - x|^2 \le (U - x, -\lambda \tilde{a}(t + \lambda h, x)h - \lambda \sum_{r=1}^m \sigma_r(t + \lambda h, x)\zeta_{rh}\sqrt{h}) \\ &\le \lambda^2 |\bar{X} - x| \left( h |\tilde{a}(t + \lambda h, x)| + \sqrt{2lh} |\ln h| \sum_{r=1}^m |\sigma_r(t + \lambda h, x)| \right). \end{split}$$

Then, using (4.3) and (4.8), we obtain (4.18), which completes the proof of Lemma 4.1 for the implicit method (4.11).

Now we consider boundedness of moments of (4.11).

LEMMA 4.2. Let  $1/2 < \lambda \leq 1$ . Assume that Assumption 4.1 holds. Then for all  $0 < h \leq h_0$  with  $h_0$  from (4.17) and for all k = 0, ..., N the following inequality holds for the fully implicit scheme (4.11) for  $p \geq 1$ :

(4.24) 
$$\mathbb{E}|X_k|^{2p} \le K(1 + \mathbb{E}|X_0|^{2p}),$$

where K > 0 is a constant.

*Proof.* We note that (4.6) and (4.9) imply

(4.25) 
$$(x, \tilde{a}(t, x)) \le (c + 3m\lambda L/2)(1 + |x|^2), t \in [t_0, T], x \in \mathbb{R},$$

which together with (4.4) ensures that the solution of (1.1) has all moments (2.5),  $p \ge 1$  [5].

Let  $U_{k+1} = (1 - \lambda)X_k + \lambda X_{k+1}$  (cf. (4.15)). We have (4.26)

$$\begin{split} V_{k+1} &:= |X_{k+1}|^2 - |X_k|^2 = 2(U_{k+1}, X_{k+1} - X_k) - (2\lambda - 1)|X_{k+1} - X_k|^2 \\ &= 2\lambda h(U_{k+1}, \tilde{a}(t_{k+\lambda}, U_{k+1})) + 2\lambda \sqrt{h} \left( U_{k+1}, \sum_{r=1}^m \sigma_r(t_{k+\lambda}, U_{k+1}) \left(\zeta_{rh}\right)_k \right) \\ &- (2\lambda - 1)h^2 |\tilde{a}(t_{k+\lambda}, U_{k+1})|^2 - (2\lambda - 1)h \left| \sum_{r=1}^m \sigma_r(t_{k+\lambda}, U_{k+1}) \left(\zeta_{rh}\right)_k \right|^2 \\ &- 2(2\lambda - 1)h^{3/2} \left( \tilde{a}(t_{k+\lambda}, U_{k+1}), \sum_{r=1}^m \sigma_r(t_{k+\lambda}, U_{k+1}) \left(\zeta_{rh}\right)_k \right) \\ &= 2\lambda (U_{k+1}, h\tilde{a}(t_{k+\lambda}, U_{k+1})) + 2\lambda \sqrt{h} \left( X_k, \sum_{r=1}^m \sigma_r(t_{k+\lambda}, U_{k+1}) \left(\zeta_{rh}\right)_k \right) \\ &+ (2\lambda^2 - 2\lambda + 1)h \left| \sum_{r=1}^m \sigma_r(t_{k+\lambda}, U_{k+1}) \left(\zeta_{rh}\right)_k \right|^2 - (2\lambda - 1)h^2 |\tilde{a}(t_{k+\lambda}, U_{k+1})|^2 \\ &+ 2(1 - \lambda)^2 h^{3/2} \left( \tilde{a}(t_{k+\lambda}, U_{k+1}), \sum_{r=1}^m \sigma_r(t_{k+\lambda}, U_{k+1}) \left(\zeta_{rh}\right)_k \right). \end{split}$$

Expanding  $\sigma_r(t_{k+\lambda}, U_{k+1})$  at  $(t_{k+\lambda}, X_k)$ , we obtain (4.27)

$$\begin{split} V_{k+1} &= 2\lambda(U_{k+1}, h\tilde{a}(t_{k+\lambda}, U_{k+1})) + 2\lambda\sqrt{h} \Biggl( X_k, \sum_{r=1}^m \sigma_r(t_{k+\lambda}, X_k) \left(\zeta_{rh}\right)_k \Biggr) \\ &+ 2\lambda\sqrt{h} \Biggl( X_k, \sum_{r=1}^m \nabla \sigma_r(t_{k+\lambda}, \theta) (U_{k+1} - X_k) \left(\zeta_{rh}\right)_k \Biggr) \\ &+ (2\lambda^2 - 2\lambda + 1)h \Biggl| \sum_{r=1}^m \sigma_r(t_{k+\lambda}, U_{k+1}) \left(\zeta_{rh}\right)_k \Biggr|^2 - (2\lambda - 1)h^2 |\tilde{a}(t_{k+\lambda}, U_{k+1})|^2 \\ &+ 2(\lambda - 1)^2 h^{3/2} \Biggl( \tilde{a}(t_{k+\lambda}, U_{k+1}), \sum_{r=1}^m \sigma_r(t_{k+\lambda}, U_{k+1}) \left(\zeta_{rh}\right)_k \Biggr) \\ &= 2\lambda(U_{k+1}, h\tilde{a}(t_{k+\lambda}, U_{k+1})) + 2\lambda\sqrt{h} \Biggl( X_k, \sum_{r=1}^m \sigma_r(t_{k+\lambda}, X_k) \left(\zeta_{rh}\right)_k \Biggr) \\ &+ 2\lambda^2 h^{3/2} \Biggl( X_k, \sum_{r=1}^m \nabla \sigma_r(t_{k+\lambda}, \theta) \tilde{a}(t_{k+\lambda}, U_{k+1}) \left(\zeta_{rh}\right)_k \Biggr) \\ &+ 2\lambda^2 h \Biggl( X_k, \sum_{r=1}^m \nabla \sigma_r(t_{k+\lambda}, \theta) \sum_{l=1}^m \sigma_l(t_{k+\lambda}, X_k) \left(\zeta_{lh}\right)_k \left(\zeta_{rh}\right)_k \Biggr) \\ &+ (2\lambda^2 - 2\lambda + 1)h \Biggl| \sum_{r=1}^m \sigma_r(t_{k+\lambda}, U_{k+1}) \left(\zeta_{rh}\right)_k \Biggr|^2 - (2\lambda - 1)h^2 |\tilde{a}(t_{k+\lambda}, U_{k+1})|^2 \\ &+ 2(1 - \lambda)^2 h^{3/2} \Biggl( \tilde{a}(t_{k+\lambda}, U_{k+1}), \sum_{r=1}^m \sigma_r(t_{k+\lambda}, U_{k+1}) \left(\zeta_{rh}\right)_k \Biggr), \end{split}$$

where  $\theta = \nu U_{k+1} - (1-\nu)X_k$ ,  $\nu \in [0,1]$ , is an intermediate point. Using (4.25), (4.10), Young's inequality, (4.4), and (4.19), we obtain

$$\begin{aligned} (4.28) \\ V_{k+1} &\leq \lambda h (2c+3\lambda mL)(1+|U_{k+1}|^2) + 2\lambda\sqrt{h} \left( X_k, \sum_{r=1}^m \sigma_r(t_{k+\lambda}, X_k) (\zeta_{rh})_k \right) \\ &+ \frac{2\lambda - 1}{2} h^2 |\tilde{a}(t_{k+\lambda}, U_{k+1})|^2 + \frac{2\lambda^4}{2\lambda - 1} h |X_k|^2 m \sum_{r=1}^m |\nabla \sigma_r(t_{k+\lambda}, \theta)|^2 |(\zeta_{rh})_k|^2 \\ &+ \lambda^2 h m |X_k|^2 \sum_{r=1}^m |\nabla \sigma_r(t_{k+\lambda}, \theta)|^2 |(\zeta_{rh})_k|^2 + \lambda^2 h m \sum_{l=1}^m |\sigma_l(t_{k+\lambda}, X_k)|^2 |(\zeta_{lh})_k|^2 \\ &+ (2\lambda^2 - 2\lambda + 1) h m \sum_{r=1}^m |\sigma_r(t_{k+\lambda}, U_{k+1})|^2 |(\zeta_{rh})_k|^2 \\ &- (2\lambda - 1) h^2 |\tilde{a}(t_{k+\lambda}, U_{k+1})|^2 \\ &+ \frac{2\lambda - 1}{2} h^2 |\tilde{a}(t_{k+\lambda}, U_{k+1})|^2 + \frac{2(1 - \lambda)^4}{2\lambda - 1} h m \sum_{r=1}^m |\sigma_r(t_{k+\lambda}, U_{k+1})|^2 |(\zeta_{rh})_k|^2 \\ &\leq \lambda h (2c+3\lambda mL)(1+|U_{k+1}|^2) + 2\lambda\sqrt{h} \left( X_k, \sum_{r=1}^m \sigma_r(t_{k+\lambda}, X_k) (\zeta_{rh})_k \right) \\ &+ \lambda^2 \left[ \frac{2\lambda^2}{2\lambda - 1} + 1 \right] h L_1^2 |X_k|^2 m \sum_{r=1}^m |(\zeta_{rh})_k|^2 + \lambda^2 h m L (1+|X_k|^2) \sum_{l=1}^m |(\zeta_{lh})_k|^2 \\ &+ \left[ 2\lambda^2 - 2\lambda + 1 + \frac{2(1 - \lambda)^4}{2\lambda - 1} \right] h m L (1+|U_{k+1}|^2) \sum_{l=1}^m |(\zeta_{lh})_k|^2. \end{aligned}$$

Then using (4.23), we arrive at

(4.29)  
$$V_{k+1} \le Kh(1+|X_k|^2) \left(1+\sum_{r=1}^m |(\zeta_{rh})_k|^2\right) + 2\lambda\sqrt{h} \left(X_k, \sum_{r=1}^m \sigma_r(t_{k+\lambda}, X_k)(\zeta_{rh})_k\right),$$

where K > 0 is independent of h and k while it depends on  $\lambda$  and on constants appearing in (4.2)–(4.10). Using (4.29), we get for integer  $p \ge 1$ 

$$(1+|X_{k+1}|^2)^p = (1+|X_k|^2+V_{k+1})^p$$

$$\leq (1+|X_k|^2)^p + K (1+|X_k|^2)^p \sum_{l=1}^p h^l \left[1+\sum_{r=1}^m |(\zeta_{rh})_k|^2\right]^l$$

$$+ p (1+|X_k|^2)^{p-1} 2\lambda h^{1/2} \left(X_k, \sum_{r=1}^m \sigma_r(t_{k+\lambda}, X_k) (\zeta_{rh})_k\right)$$

$$+ K \sum_{l=2}^p (1+|X_k|^2)^{p-l} h^{l/2} \left| \left(X_k, \sum_{r=1}^m \sigma_r(t_{k+\lambda}, X_k) (\zeta_{rh})_k\right) \right|^l,$$

whence, observing that  $X_k$  are  $\mathcal{F}_{t_k}$ -measurable while  $(\zeta_{rh})_k$  are independent of  $\mathcal{F}_{t_k}$ , it is not difficult to obtain

(

$$\mathbb{E}\left(1+|X_{k+1}|^{2}\right)^{p} \leq \mathbb{E}\left(1+|X_{k}|^{2}\right)^{p}+Kh\mathbb{E}\left(1+|X_{k}|^{2}\right)^{p} + K\sum_{l=2}^{p}\mathbb{E}\left(1+|X_{k}|^{2}\right)^{p-l}h^{l/2}\left|\left(X_{k},\sum_{r=1}^{m}\sigma_{r}(t_{k+\lambda},X_{k})\left(\zeta_{rh}\right)_{k}\right)\right|^{l} \\ \leq \mathbb{E}\left(1+|X_{k}|^{2}\right)^{p}+Kh\mathbb{E}\left(1+|X_{k}|^{2}\right)^{p},$$

which together with Gronwall's inequality completes the proof of the lemma for integer  $p \ge 1$ , and then by Jensen's inequality for noninteger p > 1 as well.

We have not succeeded in proving boundedness of moments for the midpoint scheme, i.e., (4.11) with  $\lambda = 1/2$  under Assumption 4.1. One can observe that the proof of Lemma 4.2 is not applicable to this choice of  $\lambda$  as the estimate in (4.28) blows up when  $\lambda \to 1/2$  and it is clear that the midpoint scheme is the boundary case. We also know [4] that for  $\sigma_r = 0$  (4.11) is B-stable for  $\lambda \ge 1/2$  and not B-stable (in fact, not A-stable) for  $\lambda < 1/2$ . It is natural to expect that for  $\lambda < 1/2$  the moments of (4.11) are not bounded and hence the method with  $\lambda < 1/2$  is divergent under Assumption 4.1. (See also such a conclusion for the drift-implicit  $\theta$ -method in [13].) In our experiments (section 5) the midpoint method produced accurate results.

At the same time, we proved boundedness of moments for the midpoint scheme if in addition to Assumption 4.1 we require that the diffusion coefficients  $\sigma_r(t, x)$  are bounded. The proof is similar to the proof of Lemma 3.1 and is omitted here.

LEMMA 4.3. Let the assumptions of Lemma 4.2 hold and in addition assume that the diffusion coefficients  $\sigma_r(t,x)$  are uniformly bounded. Then the moments of the midpoint method (4.11) with  $\lambda = 1/2$  have bounded moments: for  $p \geq 1$ ,

(4.30) 
$$\mathbb{E}|X_k|^{2p} \le K(1 + \mathbb{E}|X_0|^{4(p+1)\varkappa - 4})^{1/2},$$

where K > 0 is a constant.

The next lemma gives estimates for the one-step error of the method (4.11).

LEMMA 4.4. Let  $0 \leq \lambda \leq 1$ . Assume that (2.5) holds. Assume that the coefficient a(t, x) has continuous first-order partial derivative in t and in  $x^i$  and that the derivatives and the coefficient satisfy inequalities of the form (4.3); the functions  $\sigma_r(t, x)$  have continuous first-order partial derivatives in t and that the derivatives and the coefficients satisfy inequalities of the form (4.4)–(4.5); and the functions  $\nabla \sigma_r(t, x) \sigma_r(t, x)$  have continuous first partial derivatives in t and in  $x^i$  which satisfy inequalities of the form (4.5). Then the method (4.11) satisfies the inequalities (2.9) and (2.10) with  $q_1 = 2$  and  $q_2 = 1$ , respectively.

Proof of this lemma is rather routine and is similar to the global Lipschitz case [20, 21] and so is omitted here. Using Lemmas 4.1–4.4, the next proposition follows from Theorem 2.1.

PROPOSITION 4.5. Let for  $1/2 < \lambda \leq 1$  the assumptions of Lemmas 4.2 and 4.4 hold, and for  $\lambda = 1/2$  in addition assume that the diffusion coefficients  $\sigma_r(t, x)$  are uniformly bounded. Then the fully implicit method (4.11) has mean-square order 1/2, i.e., for it the inequality (2.13) holds with q = 1/2.

Remark 4.2. Consider the commutative case, i.e., when  $\Lambda_i \sigma_r = \Lambda_r \sigma_i$  (here the operator  $\Lambda_r := (\sigma_r, \partial/\partial x)$ ) or in the case of a system with one noise (i.e., m = 1). Then in the setting of Lemma 4.4, the midpoint method, i.e., (4.11) with  $\lambda = 1/2$ , satisfies the inequalities (2.9) and (2.10) with  $q_1 = 2$  and  $q_2 = 3/2$ , respectively. (See such a result in the global Lipschitz case in [20, 21].) Therefore, it converges in this case with mean-square order 1 when its moments are bounded.

5. Numerical examples. In this section we will test the following schemes: the balanced method (3.1) from section 3; the drift-implicit scheme (2.17); the fully implicit Euler scheme (4.11) with  $\lambda = 1$ ; the midpoint method (4.11) with  $\lambda = 1/2$ ; the drift-tamed Euler scheme (a modified balanced method) [10],

(5.1) 
$$X_{k+1} = X_k + h \frac{a(X_k)}{1+h |a(X_k)|} + \sum_{r=1}^m \sigma_r(t_k, X_k) \xi_{rk} \sqrt{h};$$

the fully tamed scheme [9],

(5.2) 
$$X_{k+1} = X_k + \frac{a(X_k)h + \sum_{r=1}^m \sigma_r(t_k, X_k)\xi_{rk}\sqrt{h}}{\max\left(1, h\left|ha(X_k) + \sum_{r=1}^m \sigma_r(t_k, X_k)\xi_{rk}\sqrt{h}\right|\right)};$$

and the trapezoidal scheme [21, p. 30],

(5.3) 
$$X_{k+1} = X_k + \frac{h}{2} \left[ a(X_{k+1}) + a(X_k) \right] + \sum_{r=1}^m \sigma_r(t_k, X_k) \xi_{rk} \sqrt{h}.$$

As before,  $\xi_{rk} = (w_r(t_{k+1}) - w_r(t_k))/\sqrt{h}$  are Gaussian  $\mathcal{N}(0, 1)$  i.i.d. random variables. Strong convergence with order 1/2 of (5.1) under Assumption 4.1 is proved in [10]. Strong convergence of (5.2) without order under Assumption 2.1 is proved in [9].

In all the experiments with fully implicit schemes, where the truncated random variables  $\zeta$  are used, we took l = 2 (see (4.13)). The experiments were performed using MATLAB R2012a on a Macintosh desktop computer with an Intel Xeon CPU E5462 (quad-core, 2.80 GHz). In simulations we used the Mersenne twister random generator with seed 100. Newton's method was used to solve the nonlinear algebraic equations at each step of the implicit schemes.

We test the methods on two model problems. The first one satisfies Assumption 4.1 (nonglobal Lipschitz drift, global Lipschitz diffusion) and has two noncommutative noises. The second example satisfies Assumption 2.1 (nonglobal Lipschitz both drift and diffusion). The aim of the tests is to compare performance of the methods: their accuracy (i.e., roughly speaking, size of prefactors at a power of h) and computational costs. We note that experiments cannot prove or disprove boundedness of moments of the schemes since experiments rely on a finite sample of trajectories run over a finite time interval, while blow-up of moments in divergent methods (e.g., explicit Euler scheme) is, in general, a result of large deviations [15, 22].

*Example 5.1.* Our first test model is the Stratonovich SDE of the form

(5.4) 
$$dX = (1 - X^5) dt + X \circ dw_1 + dw_2, \quad X(0) = 0.$$

In Ito's sense, the drift of the equation becomes  $a(t,x) = 1 - x^5 + x/2$ . Here we tested the balanced method (3.1), the drift-tamed scheme (5.1), the fully implicit Euler scheme (4.11) with  $\lambda = 1$ , and the midpoint method (4.11) with  $\lambda = 1/2$ . We note that for all the methods tested on this example except the midpoint rule mean-square convergence with order 1/2 is proved either in earlier papers [10, 9, 14] or here as described before.

To compute the mean-square error, we run M independent trajectories  $X^{(i)}(t)$ ,  $X_k^{(i)}$ :

(5.5) 
$$\left(E\left[X(T) - X_N\right]^2\right)^{1/2} \doteq \left(\frac{1}{M}\sum_{i=1}^M [X^{(i)}(T) - X_N^{(i)}]^2\right)^{1/2}$$

Table 5.1

Example 5.1. Mean-square errors of the selected schemes. See further details in the text.

h	$(4.11), \lambda = 1$	Rate	$(4.11), \lambda = 1/2$	Rate	(5.1)	Rate	(3.1)	Rate
0.1	1.712e-01	-	1.443e-01	_	$3.748\mathrm{e}{\text{-}01}$	-	3.594e-01	-
0.05	1.234e-01	0.47	9.224e-02	0.65	2.103e-01	0.83	3.017e-01	0.25
0.02	7.692e-02	0.52	5.261e-02	0.61	9.472e-02	0.87	2.297e-01	0.30
0.01	5.478e-02	0.49	3.549e-02	0.57	6.104 e- 02	0.63	1.778e-01	0.37
0.005	3.935e-02	0.48	2.487e-02	0.51	$3.959\mathrm{e}{\text{-}02}$	0.62	1.354e-01	0.39

We took time T = 50 and  $M = 10^4$ . The reference solution was computed by the midpoint method with small time step  $h = 10^{-4}$ . It was verified that using a different implicit scheme for simulating a reference solution does not affect the outcome of the tests. We chose the mid-point scheme as a reference since in all the experiments it produced the most accurate results.

Table 5.1 gives the mean-square errors and experimentally observed convergence rates for the corresponding methods. We checked that the number of trajectories  $M = 10^4$  was sufficiently large for the statistical errors not to significantly hinder the mean-square errors. (The Monte Carlo error computed with 95% confidence was at least 10 times smaller than the reported mean-square errors except values for (5.1) at h = 0.1 and 0.05, where it was at least 5 times smaller than the meansquare errors.) In addition to the data in the table, we evaluated errors for (3.1)for smaller time steps: h = 0.002—the error is 9.27e-02 (rate 0.41); 0.001—the error is 6.86e-02 (0.44). The observed rates of convergence of all the tested methods are close to the predicted 1/2. For a fixed time step h, the most accurate scheme is the midpoint one, and the less accurate scheme is the new balanced method (3.1). To produce the result with accuracy  $\sim 0.06 - 0.07$ , in our experiment of running  $M = 10^4$  trajectories the scheme (5.1) required 170 sec, the midpoint (4.11) with  $\lambda = 1/2$  required 329 sec, (4.11) with  $\lambda = 1$  required 723 sec, and (3.1) required 1870 sec. That is, our experiments confirmed the conclusion of [10] that the drift-tamed (modified balance) method (5.1) from [10] is highly competitive. We note that (5.1)is not applicable when diffusion grows faster than a linear function and that in this case the balanced method (3.1) can outcompete implicit schemes, as shown in the next example.

Example 5.2. Consider the SDE in the Stratonovich sense:

(5.6) 
$$dX = (1 - X^5) dt + X^2 \circ dw, \quad X(0) = 0.$$

In Ito's sense, the drift of the equation becomes  $a(t, x) = 1 - x^5 + x^3$ .

Here we tested the balanced method (3.1), the fully tamed Euler scheme (5.2), the drift-implicit scheme (2.17), the fully implicit Euler scheme (4.11) with  $\lambda = 1$ , the midpoint method (4.11) with  $\lambda = 1/2$ , and the trapezoidal scheme (5.3). We recall that in the case of nonglobal Lipschitz drift and diffusion, for the drift-implicit scheme (2.17) and the balanced method (3.1) mean-square convergence with order 1/2is shown earlier in this paper and for (2.17) it is also shown in [14]; it is not difficult to generalize the results of [13] to show boundedness of higher moments of the trapezoidal scheme (5.3) and then, using Theorem 2.1, to prove its mean-square convergence with order 1/2 (see also [14]), which is supported by the experiments. Strong convergence of (5.2) without order is proved in [9]. We note that it can be proved directly that implicit algebraic equations arising from application of the midpoint and fully implicit Euler schemes to (5.6) have unique solutions under a sufficiently small time step.



FIG. 5.1. Example 5.2. Trajectories of the fully tamed scheme (5.2) and the balanced scheme (3.1) for h = 0.1. The reference trajectory is simulated by the midpoint scheme (see (4.11) with  $\lambda = 1/2$ ) using a small time step.

The reference solution was computed by the midpoint method with small time step  $h = 10^{-4}$ . The time T = 50 and  $M = 10^4$  in (5.5).

The fully tamed scheme (5.2) did not produce accurate results until the time step size was at least h = 0.005, and we do not report its errors here, but see the remark below.

*Remark* 5.1. In contrast to the explicit balance scheme (3.1), the nature of the explicit fully tamed scheme (5.2) can lead to spurious oscillations, which significantly reduces its practical usefulness. Indeed, if at a step  $k_*$ , the event O := $|ha(X_k) + \sum_{r=1}^m \sigma_r(t_k, X_k) \xi_{rk} \sqrt{h}| > 1/h$  happens, then in the case of (5.6) the trajectory  $X_k$ ,  $k > k_*$ , oscillates approximately between  $X_{k_*}$  and  $X_{k_*} - sgn(X_{k_*})/h$ . Since the probability of the event O is positive for any step size h > 0 and grows with integration time, it is unavoidable that in some scenarios (i.e., on some trajectories) such oscillatory behavior will appear. For instance, in this experiment for h = 0.1 we observed 305 out of 1000 paths for which O happened over the time interval [0, 5]. 582 for over [0, 10], and 989 for over [0, 50]; for h = 0.05 we observed 866 out of 1000 paths over the time interval [0, 50]. Typical trajectories of the balance scheme (3.1)and the fully tamed scheme (5.2) are presented in Figure 5.1, where the reference solution is computed by the midpoint scheme with a small time step h = 0.0001. From the practical point of view, (5.2) works as long as the explicit Euler scheme works (cf. [15] and also [21, p. 17]). The strong convergence (without order) of (5.2) [9] in comparison with the explicit Euler scheme is due to the following fact. When event O happens for the Euler scheme its sequence  $X_k$  starts oscillating with growing amplitude, which leads to unboundedness of its moments and, consequently, its divergence in the mean-square sense. For (5.2), the oscillations are bounded by  $\sim 1/h$ and since the probability of O over a finite time interval is rapidly decreasing with a decrease of h, then the moments are bounded uniformly in h. At the same time, the one-step approximation of (5.2) does not satisfy the conditions (2.9) and (2.10)of Theorem 2.1. We note that the explicit balanced-type scheme (3.1) does not have such drawbacks as (5.2).

Table 5.2 gives the mean-square errors and experimentally observed convergence rates for the corresponding methods. We checked that the number of trajectories Table 5.2

Example 5.2. Mean-square errors of the selected schemes. See further details in the text.

h	(2.17)	Rate	$(4.11), \lambda = 1$	Rate	$(4.11), \lambda = 1/2$	Rate	(5.3)	Rate	(3.1)	Rate
0.2	3.449e-01	-	1.816e-01	-	1.378e-01	-	4.920e-01	-	2.102e-01	_
0.1	2.441e-01	0.50	1.331e-01	0.45	8.723e-02	0.66	3.526e-01	0.48	1.637e-01	0.36
0.05	1.592e-01	0.62	9.619e-02	0.47	5.344e-02	0.71	2.230e-01	0.66	1.270e-01	0.37
0.02	8.360e-02	0.70	6.599e-02	0.41	2.242e-02	0.95	1.048e-01	0.82	9.170e-02	0.36
0.01	5.460e-02	0.61	4.919e-02	0.42	1.145e-02	0.97	5.990e-02	0.81	7.065e-02	0.38
0.005	3.682e-02	0.57	3.522e-02	0.48	5.945e-03	0.95	3.784e-02	0.66	5.393e-02	0.39

 TABLE 5.3

 Example 5.2. Computational times for the selected schemes. See further details in the text.

h	(2.17)	$(4.11), \lambda = 1$	$(4.11), \lambda = 1/2$	(5.3)	(3.1)
0.2	9.25e + 00	1.10e+01	9.33e + 00	1.20e + 01	3.98e + 00
0.1	1.77e + 01	2.17e + 01	1.80e + 01	2.30e + 01	7.49e + 00
0.05	3.42e + 01	4.26e + 01	$3.51\mathrm{e}{+01}$	4.48e + 01	1.41e+01
0.02	8.33e+01	1.04e + 02	8.69e + 01	1.10e+02	3.37e + 01
0.01	1.64e + 02	$2.05\mathrm{e}{+02}$	1.73e + 02	$2.19\mathrm{e}{+02}$	6.62e + 01
0.005	3.25e + 02	4.07e + 02	3.47e + 02	4.37e + 02	$1.32\mathrm{e}{+02}$

 $M = 10^4$  was sufficiently large for the statistical errors not to significantly hinder the mean-square errors. (The Monte Carlo error computed with 95% confidence was at least 10 times smaller than the reported mean-square errors.) In addition to the data in the table, we evaluated errors for (3.1) for smaller time steps: for h = 0.002the error is 3.70e-02 (rate 0.41), for 0.001 it is 2.73e-02 (0.44), and for 0.0005 it is 2.00e-02 (0.45), i.e., for smaller h the observed convergence rate of (3.1) becomes closer to the theoretically predicted order 1/2. Since (5.6) is with single noise, Remark 4.2 is valid here, which explains why the midpoint scheme demonstrates the first order of convergence. The other implicit schemes show the order 1/2 as expected.

Table 5.3 presents the time costs in seconds. Let us fix the tolerance level at 0.05 - 0.06. We highlight in bold the corresponding values in both tables. We see that in this example the midpoint scheme is the most efficient, which is due to its first-order convergence in the commutative case. Among methods of order 1/2, the balanced method (3.1) is the fastest and one can expect that for multidimensional SDEs the explicit scheme (3.1) can considerably outperform implicit methods. (See a similar outcome for the drift-tamed method (5.1) supported by experiments in [10]; note that (5.1), in comparison with (3.1), is, as a rule, divergent when diffusion is growing faster than a linear function on infinity.)

Appendix A. Proof of the fundamental theorem. Note that in this and the next two sections we shall use the letter K to denote various constants which are independent of h and k. The proof exploits the idea of the proof of this theorem in the global Lipschitz case [17].

Consider the error of the method  $\bar{X}_{t_0,X_0}(t_{k+1})$  at the (k+1)-step:

(A.1) 
$$\rho_{k+1} := X_{t_0,X_0}(t_{k+1}) - X_{t_0,X_0}(t_{k+1}) = X_{t_k,X(t_k)}(t_{k+1}) - X_{t_k,X_k}(t_{k+1}) = (X_{t_k,X(t_k)}(t_{k+1}) - X_{t_k,X_k}(t_{k+1})) + (X_{t_k,X_k}(t_{k+1}) - \bar{X}_{t_k,X_k}(t_{k+1})).$$

The first difference in the right-hand side of (A.1) is the error of the solution arising due to the error in the initial data at time  $t_k$ , accumulated at the kth step, which we can rewrite as

$$S_{t_k,X(t_k),X_k}(t_{k+1}) = S_{k+1} := X_{t_k,X(t_k)}(t_{k+1}) - X_{t_k,X_k}(t_{k+1})$$
$$= \rho_k + Z_{t_k,X(t_k),X_k}(t_{k+1}) = \rho_k + Z_{k+1},$$

where Z is as in (2.14). The second difference in (A.1) is the one-step error at the (k+1)-step and we denote it as  $r_{k+1}$ :

$$r_{k+1} = X_{t_k, X_k}(t_{k+1}) - \bar{X}_{t_k, X_k}(t_{k+1}).$$

Let  $p \ge 1$  be an integer. We have

(A.2) 
$$\mathbb{E}|\rho_{k+1}|^{2p} = \mathbb{E}|S_{k+1} + r_{k+1}|^{2p} \\ = \mathbb{E}[(S_{k+1}, S_{k+1}) + 2(S_{k+1}, r_{k+1}) + (r_{k+1}, r_{k+1})]^p \\ \leq \mathbb{E}|S_{k+1}|^{2p} + 2p\mathbb{E}|S_{k+1}|^{2p-2} (\rho_k + Z_{k+1}, r_{k+1}) \\ + K \sum_{l=2}^{2p} \mathbb{E}|S_{k+1}|^{2p-l} |r_{k+1}|^l.$$

Due to (2.15) of Lemma 2.2, the first term on the right-hand side of (A.2) is estimated as

(A.3) 
$$\mathbb{E} |S_{k+1}|^{2p} \leq \mathbb{E} |\rho_k|^{2p} (1+Kh).$$

Consider the second term on the right-hand side of (A.2):

(A.4) 
$$\mathbb{E} |S_{k+1}|^{2p-2} (\rho_k + Z_{k+1}, r_{k+1}) \\ = \mathbb{E} |\rho_k|^{2p-2} (\rho_k, r_{k+1}) + \mathbb{E} \left( |S_{k+1}|^{2p-2} - |\rho_k|^{2p-2} \right) (\rho_k, r_{k+1}) \\ + \mathbb{E} |S_{k+1}|^{2p-2} (Z_{k+1}, r_{k+1}).$$

Due to  $\mathcal{F}_{t_k}$ -measurability of  $\rho_k$  and due to the conditional variant of (2.9), we get for the first term on the right-hand side of (A.4)

(A.5) 
$$\mathbb{E} |\rho_k|^{2p-2} (\rho_k, r_{k+1}) \le K \mathbb{E} |\rho_k|^{2p-1} (1+|X_k|^{2\alpha})^{1/2} h^{q_1}.$$

Consider the second term on the right-hand side of (A.4) and first note that it is equal to zero for p = 1. We have for integer  $p \ge 2$ 

$$\mathbb{E}\left(\left|S_{k+1}\right|^{2p-2} - \left|\rho_{k}\right|^{2p-2}\right)\left(\rho_{k}, r_{k+1}\right) \le K\mathbb{E}\left|Z_{k+1}\right|\left|\rho_{k}\right|\left|r_{k+1}\right| \sum_{l=0}^{2p-3} |S_{k+1}|^{2p-3-l} |\rho_{k}|^{l}.$$

Further, using  $\mathcal{F}_{t_k}$ -measurability of  $\rho_k$  and the conditional variants of (2.10), (2.15), and (2.16) and the Cauchy–Bunyakovsky–Schwarz inequality (twice), we get for  $p \geq 2$ 

(A.6) 
$$\mathbb{E}\left(|S_{k+1}|^{2p-2} - |\rho_k|^{2p-2}\right)(\rho_k, r_{k+1})$$
  
$$\leq K\mathbb{E}\left|\rho_k\right|^{2p-1}(1 + |X(t_k)|^{2\varkappa - 2} + |X_k|^{2\varkappa - 2})^{1/4}h^{q_2 + 1/2}(1 + |X_k|^{2\alpha})^{1/2}.$$

Due to  $\mathcal{F}_{t_k}$ -measurability of  $\rho_k$ , the conditional variants of (2.10) and (2.16), and the Cauchy–Bunyakovsky–Schwarz inequality (twice), we obtain for the third term on the right-hand side of (A.4)

$$(A.7) \quad \mathbb{E} |S_{k+1}|^{2p-2} (Z_{k+1}, r_{k+1}) \\ \leq \mathbb{E} \left[ \mathbb{E} \left( S_{k+1} |^{4p-4} | \mathcal{F}_{t_k} \right)^{1/2} \mathbb{E} \left( |Z_{k+1}|^4 | \mathcal{F}_{t_k} \right)^{1/4} \mathbb{E} \left( |r_{k+1}|^4 | \mathcal{F}_{t_k} \right)^{1/4} \right] \\ \leq K \mathbb{E} \left| \rho_k \right|^{2p-1} (1 + |X(t_k)|^{2\varkappa - 2} + |X_k|^{2\varkappa - 2})^{1/4} h^{q_2 + 1/2} (1 + |X_k|^{4\alpha})^{1/4} \right]$$

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Due to  $\mathcal{F}_{t_k}$ -measurability of  $\rho_k$  and due to the conditional variants of (2.10) and (2.15) and the Cauchy–Bunyakovsky–Schwarz inequality, we estimate the third term on the right-hand side of (A.2):

(A.8) 
$$K \sum_{l=2}^{2p} \mathbb{E} |S_{k+1}|^{2p-l} |r_{k+1}|^l \le K \sum_{l=2}^{2p} \mathbb{E} [\mathbb{E} (|S_{k+1}|^{4p-2l} |\mathcal{F}_{t_k})^{1/2} \mathbb{E} (|r_{k+1}|^{2l} |\mathcal{F}_{t_k})^{1/2}]$$
  
 $\le K \sum_{l=2}^{2p} \mathbb{E} [|\rho_k|^{2p-l} h^{lq_2} (1+|X_k|^{2l\alpha})^{1/2}].$ 

Substituting (A.3)–(A.8) into (A.2) and recalling that  $q_1 \ge q_2 + 1/2$ , we obtain

$$\begin{split} \mathbb{E}|\rho_{k+1}|^{2p} &\leq \mathbb{E}|\rho_{k}|^{2p}(1+Kh) + K\mathbb{E}|\rho_{k}|^{2p-1} \left(1+|X_{k}|^{2\alpha}\right)^{1/2} h^{q_{2}+1/2} \\ &+ K\mathbb{E}\left|\rho_{k}\right|^{2p-1} \left(1+|X(t_{k})|^{2\varkappa-2}+|X_{k}|^{2\varkappa-2}\right)^{1/4} h^{q_{2}+1/2} (1+|X_{k}|^{2\alpha})^{1/2} \\ &+ K\mathbb{E}\left|\rho_{k}\right|^{2p-1} \left(1+|X(t_{k})|^{2\varkappa-2}+|X_{k}|^{2\varkappa-2}\right)^{1/4} h^{q_{2}+1/2} (1+|X_{k}|^{4\alpha})^{1/4} \\ &+ K\sum_{l=2}^{2p} \mathbb{E}[|\rho_{k}|^{2p-l} h^{lq_{2}} (1+|X_{k}|^{2l\alpha})^{1/2}] \\ &\leq \mathbb{E}|\rho_{k}|^{2p} (1+Kh) \\ &+ K\mathbb{E}\left|\rho_{k}\right|^{2p-1} \left(1+|X(t_{k})|^{2\varkappa-2}+|X_{k}|^{2\varkappa-2}\right)^{1/4} h^{q_{2}+1/2} (1+|X_{k}|^{2\alpha})^{1/2} \\ &+ K\sum_{l=2}^{2p} \mathbb{E}[|\rho_{k}|^{2p-l} h^{lq_{2}} (1+|X_{k}|^{2l\alpha})^{1/2}]. \end{split}$$

Then using Young's inequality and the conditions (2.5) and (2.12), we obtain

$$\mathbb{E}|\rho_{k+1}|^{2p} \le \mathbb{E}|\rho_k|^{2p} + Kh\mathbb{E}|\rho_k|^{2p} + K(1 + \mathbb{E}|X_0|^{\beta p(\varkappa - 1) + 2p\alpha\beta})h^{2p(q_2 - 1/2) + 1},$$

whence (2.13) with integer  $p \ge 1$  follows by application of Gronwall's inequality. Then by Jensen's inequality (2.13) holds for noninteger p as well.

**Appendix B. Proof of Lemma 2.2.** Lemma 2.2 is an analogue of Lemma 1.1.3 in [21].

*Proof.* Introduce the process  $S_{t,x,y}(s) = S(s) := X_{t,x}(s) - X_{t,y}(s)$  and note that Z(s) = S(s) - (x - y). We first prove (2.15). Using the Ito formula and the condition (2.2) (recall that (2.2) implies (2.5)), we obtain for  $\theta \ge 0$ 

$$\begin{split} \mathbb{E}|S(t+\theta)|^{2p} \\ &= |x-y|^{2p} + 2p \int_{t}^{t+\theta} \mathbb{E}|S|^{2p-2} \bigg[ S^{\mathsf{T}}(a(t,X_{t,x}(s)) - a(t,X_{t,y}(s))) \\ &\quad + \frac{1}{2} \sum_{r=1}^{m} |\sigma_{r}(t,X_{t,x}(s)) - \sigma_{r}(t,X_{t,y}(s))|^{2} \bigg] ds \\ &\quad + 2p(p-1) \int_{t}^{t+\theta} \mathbb{E}|S|^{2p-4} \left| S^{\mathsf{T}}(s) \sum_{r=1}^{m} [\sigma_{r}(t,X_{t,x}(s)) - \sigma_{r}(t,X_{t,y}(s))] \right|^{2} ds \\ &\leq |x-y|^{2p} + 2p \int_{t}^{t+\theta} \mathbb{E}|S|^{2p-2} \bigg[ S^{\mathsf{T}}(a(t,X_{t,x}(s)) - a(t,X_{t,y}(s))) \bigg|^{2} ds \end{split}$$

$$+ \frac{2p-1}{2} \int_{t}^{t+\theta} \sum_{r=1}^{m} |\sigma_{r}(t, X_{t,x}(s)) - \sigma_{r}(t, X_{t,y}(s))|^{2} ds$$
  
 
$$\leq |x-y|^{2p} + 2pc_{1} \int_{t}^{t+\theta} \mathbb{E} |S(s)|^{2p} ds$$

from which (2.15) follows after applying Gronwall's inequality.

Now we prove (2.16). Using the Ito formula and the condition (2.2), we obtain for  $\theta \ge 0$ 

(B.1)

$$\begin{split} \mathbb{E} \left| Z(t+\theta) \right|^{2p} &= 2p \int_{t}^{t+\theta} \mathbb{E} |Z|^{2p-2} \bigg[ Z^{\mathsf{T}} (a(t, X_{t,x}(s)) - a(t, X_{t,y}(s))) \\ &\quad + \frac{1}{2} \sum_{r=1}^{m} |\sigma_{r}(t, X_{t,x}(s)) - \sigma_{r}(t, X_{t,y}(s))|^{2} \bigg] \, ds \\ &\quad + 2p(p-1) \int_{t}^{t+\theta} \mathbb{E} |Z|^{2p-4} \left| Z^{\mathsf{T}} \sum_{r=1}^{m} [\sigma_{r}(t, X_{t,x}(s)) - \sigma_{r}(t, X_{t,y}(s))] \right|^{2} \, ds \\ &\leq 2p \int_{t}^{t+\theta} \mathbb{E} |Z|^{2p-2} (s) \bigg[ S^{\mathsf{T}} (a(t, X_{t,x}(s)) - a(t, X_{t,y}(s)))) \\ &\quad + \frac{2p-1}{2} \int_{t}^{t+\theta} \sum_{r=1}^{m} |\sigma_{r}(t, X_{t,x}(s)) - \sigma_{r}(t, X_{t,y}(s))|^{2} \bigg] \, ds \\ &\quad - 2p \int_{t}^{t+\theta} \mathbb{E} |Z|^{2p-2} (x - y, a(t, X_{t,x}(s)) - a(t, X_{t,y}(s))) ds \\ &\leq 2pc_{1} \int_{t}^{t+\theta} \mathbb{E} |Z|^{2p-2} |S|^{2} \, ds \\ &\quad - 2p \int_{t}^{t+\theta} \mathbb{E} |Z|^{2p-2} (x - y, a(t, X_{t,x}(s)) - a(t, X_{t,y}(s))) ds. \end{split}$$

Using Young's inequality, we get for the first term in the right-hand side of (B.1)

(B.2) 
$$2pc_1 \int_t^{t+\theta} \mathbb{E}|Z|^{2p-2} |S|^2 ds \le 4pc_1 \int_t^{t+\theta} \mathbb{E}|Z|^{2p-2} (|Z|^2 + |x-y|^2) ds$$
  
 $\le K \int_t^{t+\theta} \mathbb{E}|Z|^{2p} ds + K|x-y|^2 \int_t^{t+\theta} \mathbb{E}|Z|^{2p-2} ds.$ 

Consider the second term in the right-hand side of (B.1). Using Hölder's inequality (twice), (2.3), (2.15), and (2.5), we obtain

(B.3)  

$$-2p \int_{t}^{t+\theta} \mathbb{E}|Z|^{2p-2} (x-y, a(t, X_{t,x}(s)) - a(t, X_{t,y}(s))) ds$$

$$\leq 2p \int_{t}^{t+\theta} \mathbb{E}|Z|^{2p-2} |a(t, X_{t,x}(s)) - a(t, X_{t,y}(s))| |x-y| ds$$

$$\leq K |x-y| \int_{t}^{t+\theta} \left[ \mathbb{E}|Z|^{2p} \right]^{1-1/p} \left[ \mathbb{E}|a(t, X_{t,x}(s)) - a(t, X_{t,y}(s))|^{p} \right]^{1/p} ds$$

$$\leq K|x-y|\int_{t}^{t+\theta} \left[\mathbb{E}|Z|^{2p}\right]^{1-1/p} \\ \times \left(\mathbb{E}[(1+|X_{t,x}(s)|^{2\varkappa-2}+|X_{t,y}(s)|^{2\varkappa-2})^{p/2}|X_{t,x}(s)-X_{t,y}(s)|^{p}]\right)^{1/p} ds \\ \leq K|x-y|\int_{t}^{t+\theta} \left[\mathbb{E}|Z|^{2p}\right]^{1-1/p} \left(\mathbb{E}[(1+|X_{t,x}(s)|^{2\varkappa-2}+|X_{t,y}(s)|^{2\varkappa-2})^{p}]\right)^{1/2p} \\ \times \left(\mathbb{E}[|X_{t,x}(s)-X_{t,y}(s)|^{2p}]\right)^{1/2p} ds \\ \leq K|x-y|^{2} \left(1+|x|^{2\varkappa-2}+|y|^{2\varkappa-2}\right)^{1/2} \int_{t}^{t+\theta} \left[\mathbb{E}|Z|^{2p}\right]^{1-1/p} ds.$$

Substituting (B.2) and (B.3) into (B.1) and applying Hölder's inequality to  $\mathbb{E}|Z|^{2p-2}$ , we get

$$\mathbb{E} |Z(t+\theta)|^{2p} \le K \int_{t}^{t+\theta} \mathbb{E} |Z|^{2p} ds + K |x-y|^{2} (1+|x|^{2\varkappa-2}+|y|^{2\varkappa-2})^{1/2} \int_{t}^{t+\theta} \left[\mathbb{E} |Z|^{2p}\right]^{1-1/p} ds,$$

whence we obtain (2.16) for integer  $p \ge 1$  using Gronwall's inequality as, e.g., in [25, p. 360], and then by Jensen's inequality for noninteger p > 1 as well.

**Appendix C. Proof of Lemma 3.2.** As in the global Lipschitz case [19, 21], the proof of Lemma 3.2 is routine one-step error analysis.

*Proof.* We start with proving an auxiliary result. Let a function  $\varphi(t, x)$  have continuous first-order partial derivative in t and the derivative and the function satisfy inequalities of the form (2.3). For  $\alpha \geq 1$  and  $s \geq t$ , we have

$$\begin{split} & \mathbb{E} \left| \varphi(s, X_{t,x}(s)) - \varphi(t,x) \right|^{\alpha} \\ & \leq K \mathbb{E} |\varphi(s, X_{t,x}(s)) - \varphi(s,x)|^{\alpha} + K \left| \varphi(s,x) - \varphi(t,x) \right|^{\alpha} \\ & \leq K \mathbb{E} \left[ (1 + |X_{t,x}(s)|^{\varkappa - 1} + |x|^{\varkappa - 1}) |X_{t,x}(s) - x| \right]^{\alpha} + K \left| \int_{t}^{s} \frac{\partial}{\partial s} \varphi(s',x) ds' \right|^{\alpha} \\ & \leq K \mathbb{E} (1 + |X_{t,x}(s)|^{\varkappa - 1} + |x|^{\varkappa - 1})^{\alpha} \left| \int_{t}^{s} a(s', X_{t,x}(s')) ds' \right|^{\alpha} \\ & \quad + K \sum_{r=1}^{q} \mathbb{E} (1 + |X_{t,x}(s)|^{\varkappa - 1} + |x|^{\varkappa - 1})^{\alpha} \left| \int_{t}^{s} \sigma_{r}(s', X_{t,x}(s')) dw_{r}(s') \right|^{\alpha} \\ & \quad + K (1 + |x|^{\alpha \varkappa})(s - t)^{\alpha}. \end{split}$$

Then, using the Cauchy–Bunyakovsky–Schwarz inequality, (2.5), and (2.3), we get

$$(C.1) \qquad \mathbb{E} \left| \varphi(s, X_{t,x}(s)) - \varphi(t,x) \right|^{\alpha} \\ \leq K(1+|x|^{\alpha(\varkappa-1)}) \left[ \mathbb{E} \left( \int_{t}^{s} (1+|X_{t,x}(s')|^{\varkappa}) ds' \right)^{2\alpha} \right]^{1/2} \\ + K(1+|x|^{\alpha(\varkappa-1)}) \sum_{r=1}^{q} \left[ \mathbb{E} \left| \int_{t}^{s} \sigma_{r}(s', X_{t,x}(s')) dw_{r}(s') \right|^{2\alpha} \right]^{1/2} \\ + K(1+|x|^{\alpha\varkappa})(s-t)^{\alpha}.$$

By the inequality for powers of Ito integrals from [2, p. 26], we obtain that

(C.2) 
$$\mathbb{E}\left|\int_{t}^{s} \sigma_{r}(s', X_{t,x}(s')) dw_{r}(s')\right|^{2\alpha} \leq K(s-t)^{\alpha-1} \int_{t}^{s} \mathbb{E}|\sigma_{r}(s', X_{t,x}(s'))|^{2\alpha} ds'.$$

And, by the same recipe as in [2, p. 26] which exploits Hölder's inequality, it is not difficult to get

(C.3) 
$$\mathbb{E}\left[\int_{t}^{s} |1+X_{t,x}(s')|^{\varkappa} ds'\right]^{2\alpha} \leq K(s-t)^{2\alpha-1} \int_{t}^{s} \mathbb{E}|1+X_{t,x}(s')|^{2\alpha\varkappa} ds'.$$

It follows from (C.1)–(C.3), from the assumption that  $\sigma_r$  satisfy (2.3), and from (2.5) that

(C.4) 
$$\mathbb{E} |\varphi(s, X_{t,x}(s)) - \varphi(t,x)|^{\alpha} \le K(1+|x|^{2\alpha\varkappa-\alpha})[(s-t)^{\alpha/2} + (s-t)^{\alpha}],$$

which, in particular, holds for the functions a(t, x) and  $\sigma_r(t, x)$  under the conditions of the lemma.

Now consider the one-step approximation of the SDEs (1.1), which corresponds to the balanced method (3.1),

(C.5) 
$$X = x + \frac{a(t,x)h + \sum_{r=1}^{m} \sigma_r(t,x)\xi_r\sqrt{h}}{1 + h|a(t,x)| + \sqrt{h}\sum_{r=1}^{m} |\sigma_r(t,x)\xi_r|}$$

and the one-step approximation corresponding to the explicit Euler scheme:

(C.6) 
$$\tilde{X} = x + a(t,x)h + \sum_{r=1}^{m} \sigma_r(t,x)\xi_r\sqrt{h}.$$

We start with analysis of the one-step error of the Euler scheme:

$$\tilde{\rho}(t,x) := X_{t,x}(t+h) - X.$$

Using (C.4), we obtain

$$\begin{aligned} |\mathbb{E}\tilde{\rho}(t,x)| &= \left| \mathbb{E}\int_{t}^{t+h} (a(s, X_{t,x}(s)) - a(t,x)) ds \right| \leq \mathbb{E}\int_{t}^{t+h} |a(s, X_{t,x}(s)) - a(t,x)| ds \\ &\leq Kh^{3/2} (1+|x|^{2\varkappa - 1}). \end{aligned}$$

(We remark that assuming additional smoothness of a(t, x), we can get an estimate for  $\mathbb{E}\tilde{\rho}(t, x)$  of order  $O(h^2)$ , but this will not improve the result of this lemma for the balanced scheme (3.1).) Further,

(C.8) 
$$\mathbb{E}\tilde{\rho}^{2p}(t,x) \leq K\mathbb{E}\left|\int_{t}^{t+h} (a(s,X_{t,x}(s)) - a(t,x))ds\right|^{2p} + K\sum_{r=1}^{q} \mathbb{E}\left|\int_{t}^{t+h} (\sigma_{r}(s,X_{t,x}(s)) - \sigma_{r}(t,x))dw_{r}(s)\right|^{2p}$$

Using the same recipe as in (C.3) and then applying (C.4), we get for the first term in (C.8)

(C.9)  

$$\mathbb{E}\left|\int_{t}^{t+h} (a(s, X_{t,x}(s)) - a(t,x))ds\right|^{2p} \leq Kh^{2p-1} \int_{t}^{t+h} \mathbb{E}|a(s, X_{t,x}(s)) - a(t,x)|^{2p}ds$$

$$\leq Kh^{3p}(1+|x|^{4p\varkappa-2p}).$$

Using an inequality of the form (C.2) and then applying (C.4), we obtain

(C.10) 
$$\mathbb{E} \left| \int_{t}^{t+h} \left( \sigma_{r}(s, X_{t,x}(s)) - \sigma_{r}(t, x) \right) dw_{r}(s) \right|^{2p} \le Kh^{p-1} \int_{t}^{t+h} \mathbb{E} \left| \sigma_{r}(s, X_{t,x}(s)) - \sigma_{r}(t, x) \right|^{2p} ds \le Kh^{2p} (1 + |x|^{4p\varkappa - 2p}).$$

It follows from (C.8)-(C.10) that

(C.11) 
$$\mathbb{E}\tilde{\rho}^{2p}(t,x) \le Kh^{2p}(1+|x|^{4p\varkappa-2p}).$$

Now we compare the one-step approximations (C.5) of the balanced scheme and (C.6) of the Euler scheme:

(C.12) 
$$X = x + \frac{a(t,x)h + \sum_{r=1}^{m} \sigma_r(t,x)\xi_r\sqrt{h}}{1 + h|a(t,x)| + \sqrt{h}\sum_{r=1}^{m} |\sigma_r(t,x)\xi_r|} = \tilde{X} - \rho(t,x),$$

where

$$\rho(t,x) = \left(a(t,x)h + \sum_{r=1}^{m} \sigma_r(t,x)\xi_r\sqrt{h}\right) \frac{h|a(t,x)| + \sqrt{h}\sum_{r=1}^{m} |\sigma_r(t,x)\xi_r|}{1 + h|a(t,x)| + \sqrt{h}\sum_{r=1}^{m} |\sigma_r(t,x)\xi_r|}.$$

Using the equality (3.8) and the assumptions made on the coefficients (see (2.3)), we obtain

$$|\mathbb{E}\rho(t,x)| = \left| a(t,x)h\mathbb{E}\frac{h|a(t,x)| + \sqrt{h}\sum_{r=1}^{m} |\sigma_r(t,x)\xi_r|}{1 + h|a(t,x)| + \sqrt{h}\sum_{r=1}^{m} |\sigma_r(t,x)\xi_r|} \right| \le Kh^{3/2}(1 + |x|^{2\varkappa}),$$

which together with (C.12) and (C.7) implies that (C.5) satisfies (2.9) with  $q_1 = 3/2$ . Further,

$$\mathbb{E}\rho^{2p}(t,x) \le h^{2p} \mathbb{E}\left[\sqrt{h}|a(t,x)| + \sum_{r=1}^{m} |\sigma_r(t,x)\xi_r|\right]^{4p} \le Kh^{2p}(1+|x|^{4p\varkappa}),$$

which together with (C.12) and (C.11) implies that (C.5) satisfies (2.10) with  $q_2 = 1$ .

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## REFERENCES

- N. BOU-RABEE AND E. VANDEN-EIJNDEN, A patch that imparts unconditional stability to explicit integrators for Langevin-like equations, J. Comput. Phys., 231 (2012), pp. 2565–2580.
- [2] I. I. GIHMAN AND A. V. SKOROHOD, Stochastic Differential Equations, Springer, New York, 1972.
- [3] I. GYÖNGY, A note on Euler's approximations, Potential Anal., 8 (1998), pp. 205–216.
- [4] E. HAIRER AND G. WANNER, Solving Ordinary Differential Equations. II, Springer, New York, 1996.
- [5] R. Z. HAS'MINSKIĬ, Stochastic Stability of Differential Equations, Sijthoff & Noordhoff, Groningen, the Netherlands, 1980.
- [6] D. J. HIGHAM, X. MAO, AND A. M. STUART, Strong convergence of Euler-type methods for nonlinear stochastic differential equations, SIAM J. Numer. Anal., 40 (2002), pp. 1041–1063.
- [7] D. J. HIGHAM, X. MAO, AND L. SZPRUCH, Convergence, Non-negativity and Stability of a New Milstein Scheme with Applications to Finance, arXiv:1204.1647, 2012.
- [8] Y. Hu, Semi-implicit Euler-Maruyama scheme for stiff stochastic equations, in Stochastic Analysis and Related Topics, Birkhäuser, Basel, 1996, pp. 183–202.
- M. HUTZENTHALER AND A. JENTZEN, Numerical approximation of stochastic differential equations with non-globally Lipschitz continuous coefficients, Mem. AMS., to appear.
- [10] M. HUTZENTHALER, A. JENTZEN, AND P. E. KLOEDEN, Strong convergence of an explicit numerical method for SDEs with non-globally Lipschitz continuous coefficients, Ann. Appl. Probab., 22 (2012), pp. 1611–1641.
- [11] M. HUTZENTHALER, A. JENTZEN, AND P. E. KLOEDEN, Divergence of the multilevel Monte Carlo Euler method for nonlinear stochastic differential equations, Ann. Appl. Prob., 23 (2013), pp. 1913–1966.
- [12] P. E. KLOEDEN AND E. PLATEN, Numerical Solution of Stochastic Differential Equations, Springer, New York, 1992.
- [13] X. MAO AND L. SZPRUCH, Strong convergence and stability of implicit numerical methods for stochastic differential equations with non-globally Lipschitz continuous coefficients, J. Comput. Appl. Math., 238 (2013), pp. 14–28.
- [14] X. MAO AND L. SZPRUCH, Strong convergence rates for backward Euler-Maruyama method for non-linear dissipative-type stochastic differential equations with super-linear diffusion coefficients, Stochastics, 85 (2013), pp. 144–171.
- [15] J. C. MATTINGLY, A. M. STUART, AND D. J. HIGHAM, Ergodicity for SDEs and approximations: locally Lipschitz vector fields and degenerate noise, Stochastic Process. Appl., 101 (2002), pp. 185–232.
- [16] J. H. MENTINK, M. V. TRETYAKOV, A. FASOLINO, M. I. KATSNELSON, AND T. RASING, Stable and fast semi-implicit integration of the stochastic Landau-Lifshitz equation, J. Phys. Condens. Matter, 22 (2010), 176001.
- [17] G. N. MILSTEIN, A theorem on the order of convergence of mean-square approximations of solutions of systems of stochastic differential equations, Teor. Prob. Appl., 32 (1987), pp. 809–811.
- [18] G. N. MILSTEIN, Numerical Integration of Stochastic Differential Equations, Kluwer, Dordrecht, the Netherlands, 1995.
- [19] G. N. MILSTEIN, E. PLATEN, AND H. SCHURZ, Balanced implicit methods for stiff stochastic systems, SIAM J. Numer. Anal., 35 (1998), pp. 1010–1019.
- [20] G. N. MILSTEIN, Y. M. REPIN, AND M. V. TRETYAKOV, Numerical methods for stochastic systems preserving symplectic structure, SIAM J. Numer. Anal., 40 (2002), pp. 1583–1604.
- [21] G. N. MILSTEIN AND M. V. TRETYAKOV, Stochastic Numerics for Mathematical Physics, Springer, Berlin, 2004.
- [22] G. N. MILSTEIN AND M. V. TRETYAKOV, Numerical integration of stochastic differential equations with nonglobally Lipschitz coefficients, SIAM J. Numer. Anal., 43 (2005), pp. 1139–1154.
- [23] G. N. MILSTEIN AND M. V. TRETYAKOV, Computing ergodic limits for Langevin equations, Phys. D, 229 (2007), pp. 81–95.

- [24] G. N. MILSTEIN AND M. V. TRETYAKOV, Monte Carlo methods for backward equations in nonlinear filtering, Adv. in Appl. Probab., 41 (2009), pp. 63–100.
- [25] D. S. MITRINOVIĆ AND J. E. PEČARIĆ, AND A. M. FINK, Inequalities Involving Functions and Their Integrals and Derivatives, Kluwer, Dordrecht, the Netherlands, 1991.
- [26] A. NEUENKIRCH AND L. SZPRUCH, First Order Strong Approximations of Scalar SDEs with Values in a Domain, arXiv:1209.0390, 2012.
- [27] J. M. ORTEGA AND W. C. RHEINBOLDT, Iterative Solution of Nonlinear Equations in Several Variables, Classics in Appl. Math., SIAM, Philadelphia, 2000.
- [28] J. M. SANZ-SERNA AND M. P. CALVO, Numerical Hamiltonian Problems, Chapman and Hall, London, 1994.
- [29] D. TALAY, Stochastic Hamiltonian systems: Exponential convergence to the invariant measure, and discretization by the implicit Euler scheme, Markov Process. Related Fields, 8 (2002), pp. 163–198.

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