MULTI-STEP METHODS FOR SDES AND THEIR APPLICATION TO PROBLEMS WITH SMALL NOISE*

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Abstract. In this article the numerical approximation of solutions of Itô stochastic differential equations is considered, in particular for equations with a small parameter ϵ in the noise coefficient. We construct stochastic linear multi-step methods and develop the fundamental numerical analysis concerning their mean-square consistency, numerical stability in the mean-square sense and mean-square convergence. For the special case of two-step Maruyama schemes we derive conditions guaranteeing their mean-square consistency. Further, for the small noise case we obtain expansions of the local error in terms of the step-size and the small parameter ϵ . Simulation results using several explicit and implicit stochastic linear k-step schemes, k = 1, 2, illustrate the theoretical findings.

Key words. Stochastic linear multi-step method, Mean-square convergence, Mean-square numerical stability, Mean-square consistency, Small noise, Two-step Maruyama methods.

AMS subject classifications. 60H35, 65C30, 65L06, 60H10

1. Introduction. We consider Itô stochastic differential equations (SDEs) of the form

$$X(s)\Big|_{t_0}^t = \int_{t_0}^t f(X(s), s) \, \mathrm{d}s + \int_{t_0}^t G(X(s), s) \, \mathrm{d}W(s), \quad X(t_0) = X_0, \tag{1.1}$$

for $t \in \mathcal{J}$, where $\mathcal{J} = [t_0, T]$. The drift and diffusion functions are given as $f : \mathbb{R}^n \times \mathcal{J} \to \mathbb{R}^n$, $G = (g_1, \ldots, g_m) : \mathbb{R}^n \times \mathcal{J} \to \mathbb{R}^{n \times m}$, respectively. The process W is an *m*-dimensional Wiener process on the given probability space $(\Omega, \mathcal{F}, \mathbb{P})$ with a filtration $(\mathcal{F}_t)_{t \in \mathcal{J}}$, and X_0 is a given \mathcal{F}_{t_0} -measurable initial value, independent of the Wiener process and with finite second moment. We assume that there exists a path-wise unique strong solution $X(\cdot)$ of (1.1).

The aim of this article is to analyse the mean-square convergence properties of, in general, drift-implicit linear multi-step methods (LMMs) for the approximation of the solution of (1.1). An advantage LMMs have in deterministic numerics is that they require less evaluations of the right-hand side in comparison with Runge-Kutta schemes with the same order of convergence which makes them often preferable for problems with an expensive right-hand side. We recall that in the deterministic case a high order of convergence is always based on sufficient smoothness of the solution of the differential equation. In contrast, the solution of an SDE is not smooth in the ordinary sense and a high order of convergence is achievable only by including more information on the driving Wiener process, i.e. a sufficient number of multiple stochastic integrals. In the case of general systems of SDEs, some of these integrals are difficult to simulate, in several special settings, such as additive or diagonal noise, the situation is less difficult. Our interest in stochastic linear multi-step methods (SLMMs) stems from applications with small noise in circuit simulation [11, 12, 25, 30, 31], where especially the backward differentiation formulae (BDF) have proven valuable in the deterministic case. An application in the geosciences has been reported in [13].

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Stochastic two-step methods already appear in [22] (methods of order 3/2 for equations with additive noise), in [18] and in the recent book [24]. In [2, 3] two-step methods for Itô SDEs are analysed. Stochastic versions of Adams methods for order up to five have been implemented and tested for SDEs with additive noise in [11]. Consistency of SLMMs for Stratonovich SDEs has been considered in [5], in addition stochastic Adams methods have been implemented as predictor-corrector schemes and tested. A stochastic version of Adams-Bashforth methods, which involves higher order multiple Wiener integrals, is considered in ([13]). For small noise SDEs the authors considered SLMMs incorporating mixed classical stochastic integrals in [7].

The current article provides a unified treatment of the mean-square convergence analysis of a general class of SLMMs. This class generalises all deterministic multi-step schemes as well as all Itô-schemes considered in the papers cited above. We aim at understanding when and why these schemes converge in the mean-square sense. This includes understanding the relation between local and global errors. To this end we introduce the fundamental notions of mean-square consistency, stability and convergence. The fundamental results in this respect are Theorem 3.2 and 3.3.

As an application of this result, we recapture some of the properties of deterministic LMMs for SDEs with small noise, i.e. SDEs that can be written in the form

$$X(s)\Big|_{t_0}^t = \int_{t_0}^t f(X(s), s) \, \mathrm{d}s + \int_{t_0}^t \epsilon \hat{G}(X(s), s) \, \mathrm{d}W(s), \quad t \in \mathcal{J}, \quad X(t_0) = X_0, \quad (1.2)$$

where $\epsilon \ll 1$ is a small parameter, and $\hat{G} = (\hat{g}_1, \ldots, \hat{g}_m) : \mathbb{R}^n \times \mathcal{J} \to \mathbb{R}^{n \times m}$ and its derivatives are assumed to have moderate values. For this we will derive expansions of the error in terms of the small parameter ϵ and the step-size, as has been done for one-step methods in [23]. Then the small diffusion term makes it unnecessary to use high order multiple stochastic integrals and Wiener increments will be sufficient if the step-size is not too small.

In Section 2 we introduce the class of SLMMs considered and provide necessary definitions and useful facts. In Section 3 we consider the solvability of the discrete system and the boundedness of the iterates. We then establish two fundamental results, the first one concerns the relation between mean-square numerical stability of the SLMM and Dahlquist's root condition, and the second one concerns the relation between mean-square convergence, mean-square numerical stability and mean-square consistency of the SLMM. This mirrors the results and the essential role of zero-stability in the deterministic analysis of discretisation schemes, see e.g. [10, 15, 17, 19]. In Section 4 we consider two-step-Maruyama methods and give conditions for their mean-square consistency. These conditions allow determining the parameters for the stochastic part from the parameters of the deterministic scheme and reduce to those of the underlying deterministic schemes when there is no noise. We then apply the two-step-Maruyama methods to the SDE with small noise (1.2) and derive expansions of the local error in a manner similar to that in [23]. In Section 5 we provide illustrative examples. The appendix contains the proof of Theorem 3.2.

2. Definitions and preliminary results. We denote by $|\cdot|$ the Euclidian norm in \mathbb{R}^n and by $||\cdot||$ the corresponding induced matrix norm. The mean-square norm of a vector-valued square-integrable random variable $Z \in L_2(\Omega, \mathbb{R}^n)$, with \mathbb{E} the expectation with respect to \mathbb{P} , will be denoted by

$$||Z||_{L_2} := (\mathbb{E}|Z|^2)^{1/2}$$
.

We define a deterministic grid on \mathcal{J} as $t_0 < t_1 < \ldots < t_N = T$ with (for simplicity) a constant step-size h := T/N and $t_\ell = t_0 + \ell \cdot h$, $\ell = 0, \ldots, N$. We consider a stochastic linear k-step method, which for $\ell = k, \ldots, N$, takes the form

$$\sum_{j=0}^{k} \alpha_j \ X_{\ell-j} = h \sum_{j=0}^{k} \beta_j \ f(X_{\ell-j}, t_{\ell-j}) + \sum_{j=1}^{k} \Gamma_j(X_{\ell-j}, t_{\ell-j}) \ I^{t_{\ell-j}, t_{\ell-j+1}}.$$
(2.1)

We set $\alpha_0 = 1$. We require given initial values $X_0, \ldots, X_{k-1} \in L_2(\Omega, \mathbb{R}^n)$ such that X_ℓ is \mathcal{F}_{t_ℓ} -measurable for $\ell = 0 \ldots, k-1$. As in the deterministic case, usually only $X_0 = X(t_0)$ is given by the initial value problem and the values X_1, \ldots, X_{k-1} need to be computed numerically. This can be done by suitable one-step methods, where one has to be careful to achieve the desired accuracy. Every diffusion term $\Gamma_j(x,t) I^{t,t+h}$ is a finite sum of terms each containing an appropriate function \mathcal{G} of x and t multiplied by a multiple Wiener integral over [t, t+h], i.e. it takes the general form

$$\Gamma_j(x,t) \ I^{t,t+h} = \sum_{r=1}^m \mathcal{G}_j^r(x,t) \ I_r^{t,t+h} + \sum_{\substack{r_1,r_2=0\\r_1+r_2>0}}^m \mathcal{G}_j^{r_1,r_2}(x,t) \ I_{r_1,r_2}^{t,t+h} + \dots$$

A general multiple Wiener integral is given by

$$I_{r_1,r_2,\dots,r_j}^{t,t+h}(y) = \int_t^{t+h} \int_t^{s_1} \dots \int_t^{s_{j-1}} y(X(s_j),s_j) \, \mathrm{d}W_{r_1}(s_j) \dots \, \mathrm{d}W_{r_j}(s_1), \qquad (2.2)$$

where $r_i \in \{0, 1, \ldots, m\}$ and $dW_0(s) = ds$. If $y \equiv 1$ we write $I_{r_1, r_2, \ldots, r_j}^{t, t+h}$. Note that the integral $I_r^{t, t+h}$ is simply the increment $W_r(t+h) - W_r(t)$ of the scalar Wiener process W_r . The term $I^{t, t+h}$ denotes the collection of multiple Wiener integrals associated with the interval [t, t+h]. We emphasize that an explicit discretization is used for the diffusion term. For $\beta_0 = 0$, the SLMM (2.1) is explicit, otherwise it is drift-implicit. We give two examples of two-step methods.

EXAMPLE 2.1. The first is a stochastic variant of the implicit two-step backward differentiation formula (BDF) method, which we term the BDF2-Maruyama method. For $\ell = 2, \ldots, N$, it takes the form

$$X_{\ell} - \frac{4}{3}X_{\ell-1} + \frac{1}{3}X_{\ell-2} = h \frac{2}{3}f(X_{\ell}, t_{\ell}) + \sum_{r=1}^{m} g_r(X_{\ell-1}, t_{\ell-1}) I_r^{t_{\ell-1}, t_{\ell}} - \frac{1}{3}\sum_{r=1}^{m} g_r(X_{\ell-2}, t_{\ell-2}) I_r^{t_{\ell-2}, t_{\ell-1}}.$$

Here one has $\alpha_0 = 1$, $\alpha_1 = -\frac{4}{3}$, $\alpha_2 = \frac{1}{3}$, $\beta_0 = \frac{2}{3}$, $\beta_1 = \beta_2 = 0$, and $\Gamma_1(x,t) I^{t,t+h} = \sum_{r=1}^m g_r(x,t) I^{t,t+h}_r$, $\Gamma_2(x,t) I^{t,t+h} = -\frac{1}{3} \sum_{r=1}^m g_r(x,t) I^{t,t+h}_r$.

EXAMPLE 2.2. The second example is a Milstein variant of the two-step Adams-Bashforth-method (we refer to [28, 29] for possibilities to simulate the double Wiener integrals, if the SDE is not of a type that allows simplifications in their simulation):

$$X_{\ell} - X_{\ell-1} = h \left(\frac{3}{2} f(X_{\ell-1}, t_{\ell-1}) - \frac{1}{2} f(X_{\ell-2}, t_{\ell-2}) \right) + \sum_{r=1}^{m} g_r(X_{\ell-1}, t_{\ell-1}) I_r^{t_{\ell-1}, t_{\ell}} + \sum_{q,r=1}^{m} (g_q)'_x g_r(X(t_{\ell-1}), t_{\ell-1}) I_{q,r}^{t_{\ell-1}, t_{\ell}}.$$

For this method one has $\alpha_0 = 1$, $\alpha_1 = -1$, $\alpha_2 = 0$, $\beta_0 = 0$, $\beta_1 = \frac{3}{2}$, $\beta_2 = -\frac{1}{2}$ and $\Gamma_2(x,t) \ I^{t,t+h} \equiv 0$, $\Gamma_1(x,t) \ I^{t,t+h} = \sum_{r=1}^m g_r(x,t) \ I^{t,t+h}_r + \sum_{q,r=1}^m (g_q)'_x g_r(x,t) \ I^{t,t+h}_{q,r}$.

A further example of a stochastic Adams-Bashforth method, which is covered by the general form of (2.1), is studied in [13].

We will consider mean-square convergence of SLMMs in the sense discussed in Milstein and others [1, 22, 24, 30]. Note that in the literature the term *strong* convergence is sometimes used synonymously for our expression *mean-square* convergence.

DEFINITION 2.3. We call the SLMM (2.1) for the approximation of the solution of the SDE (1.1) mean-square convergent if the global error $X(t_{\ell}) - X_{\ell}$ satisfies

$$\max_{\ell=1,...,N} \|X(t_{\ell}) - X_{\ell}\|_{L_2} \to 0 \text{ as } h \to 0,$$

we say it is mean-square convergent with order γ ($\gamma > 0$) if the global error satisfies

$$\max_{\ell=1,...,N} \|X(t_{\ell}) - X_{\ell}\|_{L_{2}} \le C \cdot h^{\gamma},$$

with constant C > 0 which is independent of the step-size h.

In the following we will define what we understand by local errors. We would like to point out that for the analysis of one-step schemes essentially two different but related concepts are used in the literature. In the first one the local error is defined as the defect that is obtained when the exact solution values are inserted into the numerical scheme. In the second one the local error is defined as the difference after one step of the exact and the numerical solution started at an arbitrary deterministic value. These concepts differ in the way the error is transported to the end of the integration interval, in the first via the numerical method, in the second via the exact solution. The second definition has been used in the fundamental work of Milstein in [21, 22], where for the first time the relation between local and global errors of one-step methods for SDEs has been clarified. However, only the first definition extends easily to multi-step methods, hence we will use it here. For comparison of these principles in the deterministic setting see [16, Chapters II.3,III.4].

DEFINITION 2.4. We define the local error of the SLMM (2.1) for the approximation of the solution of the SDE (1.1), for $\ell = k, ..., N$, as

$$L_{\ell} := \sum_{j=0}^{k} \alpha_{j} X(t_{\ell-j}) - h \sum_{j=0}^{k} \beta_{j} f(X(t_{\ell-j}), t_{\ell-j}) - \sum_{j=1}^{k} \Gamma_{j}(X(t_{\ell-j}), t_{\ell-j}) I^{t_{\ell-j}, t_{\ell-j+1}}.$$
(2.3)

We aim to conclude mean square convergence from local properties of the SLMM by means of numerical stability in the mean-square sense. Numerical stability concerns the influence of perturbations of the right-hand side of the discrete scheme on the global solution of that discrete scheme. Sources of perturbations may be the local error, round-off errors or defects in the approximate solution of implicit schemes. The mean-square stability estimate of the global error is based on the mean square norm **and** on the conditional mean of the perturbations. In the case of one-step schemes this appears e.g. in [1, 30], we refer in particular to the discussion in [22, Chapter 1.4]. We remark that in the case of k-step schemes the conditional mean has to be taken with respect to the σ -algebra $\mathcal{F}_{t_{\ell-k}}$. In our analysis we thus consider the following discrete system, the perturbed form of (2.1), for $\ell = k, \ldots, N$

$$\sum_{j=0}^{k} \alpha_j \tilde{X}_{\ell-j} = h \sum_{j=0}^{k} \beta_j f(\tilde{X}_{\ell-j}, t_{\ell-j}) + \sum_{j=1}^{k} \Gamma_j(\tilde{X}_{\ell-j}, t_{\ell-j}) I^{t_{\ell-j}, t_{\ell-j+1}} + D_\ell, \quad (2.4)$$

with initial values $\tilde{X}_{\ell} = X_{\ell} + D_{\ell}, \ \ell = 0, \dots, k-1$. We suppose that the perturbations D_{ℓ} are $\mathcal{F}_{t_{\ell}}$ -measurable and that $D_{\ell} \in L_2(\Omega, \mathbb{R}^n)$.

REMARK 2.5. It is useful to represent the perturbations in the form

$$D_{\ell} = R_{\ell} + S_{\ell} =: R_{\ell} + \sum_{j=1}^{k} S_{j,\ell-j+1}, \quad \ell = k, \dots, N,$$

where each $S_{j,\ell}$ is $\mathcal{F}_{t_{\ell}}$ -measurable with $\mathbb{E}(S_{j,\ell}|\mathcal{F}_{t_{\ell-1}}) = 0.$ (2.5)

The representation (2.5) is not unique. One extreme possibility is $R_{\ell} = D_{\ell}$, $S_{\ell} = 0$, another, more useful one, is given by

$$\begin{array}{lll}
R_{\ell}^{*} &=& \mathbb{E}(D_{\ell}|\mathcal{F}_{t_{\ell-k}}), & S_{\ell}^{*} = D_{\ell} - R_{\ell}^{*}, \\
S_{j,\ell-j+1}^{*} &=& \mathbb{E}(D_{\ell} - R_{\ell}^{*} - \sum_{i=j+1}^{k} S_{i,\ell-i+1}^{*}|\mathcal{F}_{t_{\ell-j+1}}), \ j = k, k-1, \dots, 1.
\end{array}$$
(2.6)

This construction guarantees the required measurability conditions in (2.5). We also note that this decomposition is orthogonal in $L_2(\Omega)$, i.e.

$$\|D_{\ell}\|_{L_{2}}^{2} = \|R_{\ell}^{*}\|_{L_{2}}^{2} + \sum_{j=1}^{k} \|S_{j,\ell-j+1}^{*}\|_{L_{2}}^{2}.$$
(2.7)

As an example one obtains for k = 3

$$R_{\ell}^{*} = \mathbb{E}(D_{\ell}|\mathcal{F}_{t_{\ell-3}})$$

$$S_{3,\ell-2}^{*} = \mathbb{E}(D_{\ell} - R_{\ell}^{*}|\mathcal{F}_{t_{\ell-2}})$$

$$S_{2,\ell-1}^{*} = \mathbb{E}(D_{\ell} - R_{\ell}^{*} - S_{3,\ell-2}^{*}|\mathcal{F}_{t_{\ell-1}})$$

$$S_{1,\ell}^{*} = D_{\ell} - R_{\ell}^{*} - S_{3,\ell-2}^{*} - S_{2,\ell-1}^{*}.$$

Here, in the hypothetical case that $D_{\ell} = c_0 I_r^{t_{\ell-1},t_{\ell}} + c_1 I_r^{t_{\ell-2},t_{\ell-1}} + c_2 I_r^{t_{\ell-3},t_{\ell-2}} + c_3$, we have $R_{\ell}^* = c_3$, $S_{3,\ell-2}^* = c_2 I_r^{t_{\ell-3},t_{\ell-2}}$, $S_{2,\ell-1}^* = c_1 I_r^{t_{\ell-2},t_{\ell-1}}$, $S_{1,\ell}^* = c_0 I_r^{t_{\ell-1},t_{\ell}}$.

Now we give the precise definition of mean square stability and consistency that we consider in this paper.

DEFINITION 2.6. We call the SLMM (2.1) numerically stable in the meansquare sense if there exist constants $h_0 > 0$ and S > 0 such that for all stepsizes $h < h_0$ and for all \mathcal{F}_{t_ℓ} measurable perturbations $D_\ell \in L_2(\Omega, \mathbb{R}^n)$ ($\ell = 0, ..., N$) and all their representations (2.5), the following inequality holds

$$\max_{\ell=0,\dots,N} \|X_{\ell} - \tilde{X}_{\ell}\|_{L_{2}} \le S \Big\{ \max_{\ell=0,\dots,k-1} \|D_{\ell}\|_{L_{2}} + \max_{\ell=k,\dots,N} \Big(\frac{\|R_{\ell}\|_{L_{2}}}{h} + \frac{\|S_{\ell}\|_{L_{2}}}{h^{1/2}} \Big) \Big\}, \quad (2.8)$$

where $(X_{\ell})_{\ell=1}^{N}$ and $(\tilde{X}_{\ell})_{\ell=1}^{N}$ are the solutions of the SLMM (2.1) and the perturbed discrete system (2.4), respectively.

We refer to S as the stability constant and to (2.8) as the stability inequality.

DEFINITION 2.7. We call the SLMM (2.1) for the approximation of the solution of the SDE (1.1) mean-square consistent if the local error L_{ℓ} satisfies

$$h^{-1} \|\mathbb{E}(L_{\ell}|\mathcal{F}_{t_{\ell-k}})\|_{L_2} \to 0 \text{ for } h \to 0, \text{ and } h^{-1/2} \|L_{\ell}\|_{L_2} \to 0 \text{ for } h \to 0.$$

We call the SLMM (2.1) for the approximation of the solution of the SDE (1.1) mean-square consistent of order γ ($\gamma > 0$), if the local error L_{ℓ} satisfies

$$\|\mathbb{E}(L_{\ell}|\mathcal{F}_{t_{\ell-k}})\|_{L_{2}} \leq \bar{c} \cdot h^{\gamma+1}, and \|L_{\ell}\|_{L_{2}} \leq c \cdot h^{\gamma+\frac{1}{2}}, \ell = k, \dots, N,$$

with constants $c, \bar{c} > 0$ only depending on the SDE and its solution. We remind the reader that consistency is only concerned with the local error. In

the case that we disregard other sources of errors in (2.4) we only have to deal with perturbations $D_{\ell} = L_{\ell}$.

LEMMA 2.8. The SLMM (2.1) is mean-square consistent of order γ , if

$$||R_{\ell}||_{L_{2}} \leq \bar{c} \cdot h^{\gamma+1}$$
, and $||S_{\ell}||_{L_{2}} \leq c \cdot h^{\gamma+\frac{1}{2}}$, $\ell = k, \dots, N$,

for any representation (2.5) of the local error $D_{\ell} = L_{\ell}$. The SLMM (2.1) is meansquare consistent of order γ , if and only if

$$||R_{\ell}^*||_{L_2} \le \bar{c} \cdot h^{\gamma+1}$$
, and $||S_{\ell}^*||_{L_2} \le c \cdot h^{\gamma+\frac{1}{2}}$, $\ell = k, \dots, N$,

where the representation (2.6) is chosen for the local error $D_{\ell} = L_{\ell}$. Proof. First, let $L_{\ell} = R_{\ell} + S_{\ell} =: R_{\ell} + \sum_{j=1}^{k} S_{j,\ell-j+1}$ be a representation (2.5) with $\|R_{\ell}\|_{L_2} \leq \bar{c} \cdot h^{\gamma+1}$, and $\|S_{\ell}\|_{L_2} \leq c \cdot h^{\gamma+\frac{1}{2}}$, $\ell = k, \ldots, N$. By the conditions $\mathbb{E}(S_{j,\ell}|\mathcal{F}_{t_{\ell-1}}) = 0$ we conclude

$$\|\mathbb{E}(L_{\ell}|\mathcal{F}_{t_{\ell-k}})\|_{L_{2}} = \|\mathbb{E}(R_{\ell}|\mathcal{F}_{t_{\ell-k}})\|_{L_{2}} \le \|R_{\ell}\|_{L_{2}} \le \bar{c} \cdot h^{\gamma+1}.$$

Further, we have, for $h \leq 1$,

$$\|L_{\ell}\|_{L_{2}} = \|R_{\ell} + S_{\ell}\|_{L_{2}} \le \|R_{\ell}\|_{L_{2}} + \|S_{\ell}\|_{L_{2}} \le \bar{c} \cdot h^{\gamma+1} + c \cdot h^{\gamma+\frac{1}{2}} \le (\bar{c} + c) \cdot h^{\gamma+\frac{1}{2}}.$$

Second, let the SLMM (2.1) be mean-square consistent of order γ , i.e. $\|\mathbb{E}(L_{\ell}|\mathcal{F}_{t_{\ell-k}})\|_{L_2} \leq \bar{c} \cdot h^{\gamma+1}$, and $\|L_{\ell}\|_{L_2} \leq c \cdot h^{\gamma+\frac{1}{2}}$, $\ell = k, \ldots, N$. Because of $R_{\ell}^* = \mathbb{E}(L_{\ell}|\mathcal{F}_{t_{\ell-k}})$ we then, obviously, have $\|R_{\ell}^*\|_{L_2} \leq \bar{c} \cdot h^{\gamma+1}$, and, further

$$\|S_{\ell}^{*}\|_{L_{2}} = \|L_{\ell}\|_{L_{2}} - \|R_{\ell}^{*}\|_{L_{2}} \le \|L_{\ell}\|_{L_{2}} \le c \cdot h^{\gamma + \frac{1}{2}}.$$

DEFINITION 2.9. A function $f : \mathbb{R}^n \times \mathcal{J} \to \mathbb{R}^n$ satisfies a uniform Lipschitz condition with respect to x, if there exists a positive constant L_f , such that

For further reference we state the following definitions and results.

$$|f(x,t) - f(y,t)| \le L_f |x - y|, \quad \forall x, y \in \mathbb{R}^n, t \in \mathcal{J}.$$
(2.9)

A function $\Gamma : \mathbb{R}^n \times \mathcal{J} \to \mathbb{R}^{n \times m_{\Gamma}}$ satisfies a uniform Lipschitz condition with respect to x, if there exists a positive constant L_{Γ} , such that

$$\|\Gamma(x,t) - \Gamma(y,t)\| \le L_{\Gamma}|x-y| \quad \forall x,y \in \mathbb{R}^n, t \in \mathcal{J}.$$
(2.10)

Let $C^{s,s-1}$ denote the class of all functions from $\mathbb{R}^n \times \mathcal{J}$ to \mathbb{R}^n having continuous partial derivatives up to order s-1 and, in addition, continuous partial derivatives

of order s with respect to the first variable.

Let C^K denote the class of functions y from $\mathbb{R}^n \times \mathcal{J}$ to \mathbb{R}^n that satisfy a linear growth condition in the form

$$|y(x,t)| \le K(1+|x|^2)^{\frac{1}{2}}, \qquad \forall x \in \mathbb{R}^n, t \in \mathcal{J}.$$
(2.11)

DEFINITION 2.10. The characteristic polynomial of (2.1) is given by

$$\rho(\zeta) = \alpha_0 \,\zeta^k + \alpha_1 \zeta^{k-1} + \ldots + \alpha_k. \tag{2.12}$$

The SLMM (2.1) is said to fulfil Dahlquist's root condition if

i) the roots of $\rho(\zeta)$ lie on or within the unit circle;

ii) the roots on the unit circle are simple.

LEMMA 2.11. (A discrete version of Gronwall's lemma) Let a_{ℓ} , $\ell = 1, ..., N$, and C_1, C_2 be nonnegative real numbers and assume that the inequalities

$$a_{\ell} \leq C_1 + C_2 \frac{1}{N} \sum_{i=1}^{\ell-1} a_i, \quad \ell = 1, \dots, N,$$

are valid. Then we have $\max_{\ell=1,\ldots,N} a_{\ell} \leq C_1 \exp(C_2).$

To estimate the multiple integrals (2.2) we will use the following lemma (cf. Lemma 2.1 and 2.2 in [22]).

LEMMA 2.12. For any function y belonging to the class C^K , and any $t \in \mathcal{J}, h > 0$, such that $t+h \in \mathcal{J}$, we have that

$$\mathbb{E}(I_{r_1...r_i}^{t,t+h}(y)|\mathcal{F}_t) = 0 \quad if \ r_i \neq 0 \ for \ some \ i \in \{1,...,j\},$$
(2.13)

$$\mathbb{E}(I_{r_1...,r_j}^{t,t+h}(y)|\mathcal{F}_t) = 0 \quad if \; r_i \neq 0 \; \text{for some } i \in \{1,...,j\},$$

$$\mathbb{E}(I_{r_1...,r_j}^{t,t+h}(y)|\mathcal{F}_t) \|_{L_2} \leq \|I_{r_1,...,r_j}^{t,t+h}(y)\|_{L_2} = \mathcal{O}(h^{l_1+l_2/2}),$$

$$(2.14)$$

where l_1 is the number of zero indices r_i and l_2 the number of non-zero indices r_i .

3. Global properties of stochastic LMMs. In this section we will first establish the solvability of the recurrence equations (2.4) (and thus of (2.1)), then we will discuss numerical stability and mean-square convergence of the SLMM (2.1). The former characterizes the robustness of a numerical scheme with respect to small perturbations, such as rounding errors. As a property of the numerical scheme alone, it is not a priori giving evidence on the approximation power of the scheme (which may very well approximate a different problem than intended). However, numerical stability and consistency together yield convergence of the numerical solution to the exact solution. In order to distinguish this stability concept from others, it is sometimes called zero-stability or in honour of Dahlquist, also D-stability. It should not be mistaken for properties like asymptotic stability, which guarantee that for fixed step-sizes (and long or unbounded time-intervals) qualitative properties of the exact solutions like damping behaviour in dissipative systems are preserved by the discrete approximations. For further discussions we refer the reader to the deterministic literature (see e.g. [10, 19, 20]). In the stochastic literature mean-square numerical stability for one-step-schemes has been considered in [1, 3, 4, 14, 27, 30]. In the first four of these works only perturbations in the initial data have been treated.

We now turn to the solvability of the recurrence equations. If in (2.1) and (2.4) the parameter $\beta_0 = 0$, the discrete systems are explicit and every iterate X_{ℓ} , $\ell \ge k$, can

be obtained explicitly for given $I^{t_{\ell},t_{\ell+1}}$, i.e. the recurrence equations (2.1) and (2.4) obviously have unique solutions. In the case of implicit systems we need to consider the solvability of the systems of nonlinear equations (2.1) and (2.4). In addition, we have to verify that the mean-square norm of the iterates exists. (The straightforward extension to fully implicit systems would serve as an example where the mean-square norm of the iterates does not exist.)

THEOREM 3.1. Suppose that $\beta_0 \neq 0$ and the drift-coefficient f satisfies (2.9) and assume that $2h\beta_0 L_f < 1$. Then the perturbed discrete scheme (2.4) and, in consequence, the SLMM (2.1) have a unique solution. If, in addition, the coefficients Γ_j satisfy (2.10), then the mean-square norm of the iterates exists.

Proof. The proof of the existence of unique solutions of the perturbed discrete system (2.4) (and thus of (2.1)) follows the line of proofs used in the deterministic analysis of multi-step schemes. The idea is to express (2.4) as

$$\dot{X}_{\ell} = h\beta_0 f(\dot{X}_{\ell}, t_{\ell}) + \dot{B}_{\ell} ,$$
(3.1)

where

$$\tilde{B}_{\ell} := -\sum_{j=1}^{k} \alpha_j \tilde{X}_{\ell-j} + h \sum_{j=1}^{k} \beta_j f(\tilde{X}_{\ell-j}, t_{\ell-j}) + \sum_{j=1}^{k} \Gamma_j(X_{\ell-j}, t_{\ell-j}) I^{t_{\ell-j}, t_{\ell-j+1}} + \tilde{D}_{\ell}$$

is a known $\mathcal{F}_{t_{\ell}}$ -measurable random variable, when we suppose that $X_{\ell-j}$ are known $\mathcal{F}_{t_{\ell-j}}$ -measurable random variables for $j = 1, \ldots, k$. We can then view (2.4) as a fixed-point equation in x

$$x = h\beta_0 f(x, t_\ell) + b_\ell$$

and apply the contraction mapping principle. We refer e.g. to [15, Thm.6.1.1] for more details.

It remains to show that the second moments of the iterates exist. We start from the assumption $\mathbb{E}|\tilde{X}_{\ell-j}|^2 < \infty$, $j = 1, \ldots, k$ on the initial values. Recursively, we conclude that $\mathbb{E}|\tilde{X}_{\ell}|^2 < \infty$, $\ell = k, \ldots, N$ by comparing \tilde{X}_{ℓ} with the solution of the fixed-point equation (3.1) for $\tilde{B}_{\ell} = 0$, i.e., with the solution of the deterministic implicit equation $x = h\beta_0 f(x, t_{\ell})$, and applying Lipschitz continuity arguments. For more details we refer to [30, Thm.5].

We now formulate our main theorem on numerical stability. The proof is given in the appendix.

THEOREM 3.2. The stochastic linear multi-step method (2.1) is numerically stable in the mean-square sense for every continuous f and Γ_j satisfying (2.9) and (2.10), respectively, if and only if its characteristic polynomial $\rho(\zeta)$ (2.12) satisfies Dahlquist's root condition given in Definition 2.10.

With the powerful notion of numerical stability in the mean-square sense, together with mean-square consistency the mean-square convergence follows almost immediately.

THEOREM 3.3. A mean-square consistent SLMM (2.1) for the approximation of the solution of SDE (1.1) is mean-square convergent for all continuous f and Γ_j satisfying (2.9) and (2.10), respectively, if and only if it is numerically stable in the mean-square sense. If, in addition, it is mean-square consistent with order $\gamma > 0$, then the SLMM (2.1) is mean-square convergent with order γ .

Proof. First, let us assume that the mean-square consistent numerical method (2.1) is mean-square convergent. Then the necessity of stability can essentially be proved

as in the deterministic case. Set $f \equiv 0$, $\Gamma_j \equiv 0$, $X_0 = 0$. Then (2.1) reduces to $\sum_{j=1}^{k} \alpha_j X_{\ell-j} = 0$, $l = k, k+1, \ldots$, a deterministic homogeneous difference equation, and stability follows by standard arguments, see e.g. [15].

Second, let us assume that the numerical method (2.1) is mean-square stable and consistent with order $\gamma > 0$. Then mean square convergence with order γ follows by applying the stability estimate (2.8) to $\{\tilde{X}_{\ell} := X(t_{\ell})\}$ related to the perturbations $\{D_{\ell} := L_{\ell} = R_{\ell}^* + S_{\ell}^*\}$.

4. Two-step-Maruyama schemes. In this section we consider linear two-step-Maruyama schemes, thus we have for $\ell = 2, ..., N$

$$\sum_{j=0}^{2} \alpha_j X_{\ell-j} = h \sum_{j=0}^{2} \beta_j f(X_{\ell-j}, t_{\ell-j}) + \sum_{j=1}^{2} \gamma_j \sum_{r=1}^{m} g_r(X_{\ell-j}, t_{\ell-j}) \ I_r^{t_{\ell-j}, t_{\ell-j+1}}.$$
 (4.1)

For drift and diffusion coefficients f, g_1, \ldots, g_m , which are continuous and satisfy (2.9), Theorem 3.2 applies and the two-step scheme (4.1) is mean square stable if the coefficients α_0 , α_1 , α_2 satisfy Dahlquist's root condition. If, additionally, the scheme is mean-square consistent of order γ , which in general requires more smoothness of the coefficient functions, then the scheme (4.1) is mean-square convergent of that order. Thus we will be concerned with mean-square consistency of the above scheme and derive order conditions in terms of the coefficients $\alpha_0, \alpha_1, \alpha_2, \beta_0, \beta_1, \beta_2, \gamma_1, \gamma_2$. In general, the mean-square order of convergence will be only $\frac{1}{2}$, since the only information about the driving noise process that the Maruyama-type schemes include are the Wiener increments. We note that the simple Euler-Maruyama method would suffice to obtain the same order of convergence. However, convergence is an asymptotic property, i.e. it holds for $h \to 0$ and a result concerning the order of convergence may not provide sufficient information about the size of the actual error that arise for reasonable choices of the step-size. In particular when one considers equations with a small noise term as in Equation (1.2), one may find that the influence of the noise is not dominant and properties of the methods in the deterministic setting are recovered to a certain extent.

From the deterministic theory we know that for linear k-step methods

$$\sum_{j=0}^{k} \alpha_j \, x_{\ell-j} = h \sum_{j=0}^{k} \beta_j \, f(x_{\ell-j}, t_{\ell-j}), \quad \text{applied to } x'(t) = f(x(t), t),$$

the local error is of order p + 1 for sufficiently smooth functions f if

$$\sum_{j=0}^{k} \alpha_j = 0 \quad \text{and} \quad \sum_{j=0}^{k} \alpha_j (k-j)^q = q \sum_{j=0}^{k} \beta_j (k-j)^{q-1} \quad \text{for } q = 1, \dots, p.$$

In the first part of this section we derive consistency conditions for the two-step scheme (4.1) applied to the general SDE (1.1). We establish a representation of the local error L_{ℓ} in terms of certain multiple stochastic integrals obtained by the Itô-Taylor expansion. It turns out that consistency is guaranteed under the above conditions for deterministic order 1 and additional conditions that determine the method parameters γ_1 and γ_2 .

In the second part of this section we consider the application of the scheme (4.1) in case that the noise is small. The smallness of the noise is measured by means of the parameter ϵ in the diffusion coefficient $G(x,t) = \epsilon \hat{G}(x,t)$. We emphasize that the

numerical schemes include only values of G, the explicit dependence of the diffusion coefficient on the parameter ϵ is used only for a discussion of the errors. We follow ideas of [23] and develop the local error in powers of the stepsize h and the small parameter ϵ . The expansion yields the deterministic conditions for order two and we discuss for which choices of ϵ and h the stochastic component in the error estimates becomes small compared to these order 2 terms.

4.1. Two-step schemes for general SDEs. To analyse the local error L_{ℓ} of the scheme (4.1) for the SDE (1.1) and to achieve a suitable representation (2.5) we will derive appropriate Itô-Taylor expansions, where we take special care to separate the stochastic integrals over the different subintervals of integration. We introduce operators Λ_0 and Λ_r , $r = 1, \ldots, m$, defined on $C^{2,1}$ and $C^{1,0}$, respectively, by

$$\Lambda_0 y = y'_t + y'_x f + \frac{1}{2} \sum_{r=1}^m y''_{xx}[g_r, g_r] , \quad \Lambda_r y = y'_x g_r , \ r = 1, \dots, m,$$
 (4.2)

and remind the reader of the notation for multiple Wiener integrals (2.2). Using these operators the Itô formula for a function y in $C^{2,1}$ and the solution X of (1.1) reads

$$y(X(t),t) = y(X(t_0),t_0) + I_0^{t_0,t}(\Lambda_0 y) + \sum_{r=1}^m I_r^{t_0,t}(\Lambda_r y), \quad t \in \mathcal{J}.$$
 (4.3)

Applying the Itô-formula (4.3) on the corresponding intervals to the drift coefficient f as well as to the diffusion coefficients g_r yields for $s \in [t_{\ell-j}, t_{\ell-j+1}], j = 1, 2$

$$f(X(s),s) = f(X(t_{\ell-j}), t_{\ell-j}) + I_0^{t_{\ell-j},s}(\Lambda_0 f) + \sum_{r=1}^m I_r^{t_{\ell-j},s}(\Lambda_r f),$$
(4.4)

$$g_r(X(s),s) = g_r(X(t_{\ell-j}), t_{\ell-j}) + I_0^{t_{\ell-j},s}(\Lambda_0 g_r) + \sum_{q=1}^m I_q^{t_{\ell-j},s}(\Lambda_q g_r).$$
(4.5)

We trace back the values of the drift coefficient to the point $t_{\ell-2}$ to obtain

$$f(X(t_{\ell-1}), t_{\ell-1}) = f(X(t_{\ell-2}), t_{\ell-2}) + I_0^{t_{\ell-2}, t_{\ell-1}}(\Lambda_0 f) + \sum_{r=1}^m I_r^{t_{\ell-2}, t_{\ell-1}}(\Lambda_r f), (4.6)$$

$$f(X(t_{\ell}), t_{\ell}) = f(X(t_{\ell-2}), t_{\ell-2}) + I_0^{t_{\ell-2}, t_{\ell-1}}(\Lambda_0 f) + I_0^{t_{\ell-1}, t_{\ell}}(\Lambda_0 f)$$

$$+ \sum_{r=1}^m I_r^{t_{\ell-2}, t_{\ell-1}}(\Lambda_r f) + \sum_{r=1}^m I_r^{t_{\ell-1}, t_{\ell}}(\Lambda_r f).$$

$$(4.7)$$

For the general SDE (1.1) we have the following result.

LEMMA 4.1. Assume that the coefficients $f, g_r, r = 1, ..., m$ of the SDE (1.1) belong to the class $C^{2,1}$ with $\Lambda_0 f, \Lambda_0 g_r, \Lambda_r f, \Lambda_q g_r \in C^K$ for r, q = 1, ..., m. Then the local error (2.3) of the stochastic 2-step scheme (4.1) allows the representation

$$L_{\ell} = R_{\ell}^{\circ} + S_{1,\ell}^{\circ} + S_{2,\ell-1}^{\circ}, \quad \ell = 2, \dots, N,$$
(4.8)

where $R_{\ell}^{\circ}, S_{j,\ell}^{\circ}, j = 1, 2$ are $\mathcal{F}_{t_{\ell}}$ -measurable with $\mathbb{E}(S_{j,\ell}^{\circ}|\mathcal{F}_{t_{\ell-1}}) = 0$ and

$$R_{\ell}^{\circ} = \left[\sum_{j=0}^{2} \alpha_{j}\right] X(t_{\ell-2}) + \left[2\alpha_{0} + \alpha_{1} - \sum_{j=0}^{2} \beta_{j}\right] hf(X(t_{\ell-2}), t_{\ell-2}) + \tilde{R}_{\ell}^{\circ},$$

$$S_{1,\ell}^{\circ} = \left[\alpha_{0} - \gamma_{1}\right] \sum_{r=1}^{m} g_{r}(X(t_{\ell-1}), t_{\ell-1}) I_{r}^{t_{\ell-1}, t_{\ell}} + \tilde{S}_{1,\ell}^{\circ},$$

$$S_{2,\ell-1}^{\circ} = \left[(\alpha_{0} + \alpha_{1}) - \gamma_{2}\right] \sum_{r=1}^{m} g_{r}(X(t_{\ell-2}), t_{\ell-2}) I_{r}^{t_{\ell-2}, t_{\ell-1}} + \tilde{S}_{2,\ell-1}^{\circ}$$

with

$$\|\tilde{R}^{\circ}_{\ell}\|_{L_{2}} = O(h^{2}), \quad \|\tilde{S}^{\circ}_{1,\ell}\|_{L_{2}} = O(h), \quad \|\tilde{S}^{\circ}_{2,\ell-1}\|_{L_{2}} = O(h).$$
(4.9)

COROLLARY 4.2. Let the coefficients f, g_r , r = 1, ..., m, of the SDE (1.1) satisfy the assumptions of Lemma 4.1 and suppose they are Lipschitz continuous with respect to their first variable. Let the coefficients of the stochastic linear two-step scheme (4.1) satisfy Dahlquist's root condition and the consistency conditions

$$\sum_{j=0}^{2} \alpha_j = 0, \quad 2\alpha_0 + \alpha_1 = \sum_{j=0}^{2} \beta_j, \quad \alpha_0 = \gamma_1, \quad \alpha_0 + \alpha_1 = \gamma_2.$$
(4.10)

Then the global error of the scheme (4.1) applied to (1.1) allows the expansion

$$\max_{\ell=2,\dots,N} \|X(t_{\ell}) - X_{\ell}\|_{L_{2}} = \mathcal{O}(h^{1/2}) + \mathcal{O}(\max_{\ell=0,1} \|X(t_{\ell}) - X_{\ell}\|_{L_{2}}) .$$

Proof. (of Corollary 4.2) By Lemma 4.1 we have the representation (4.8) for the local error. Applying the consistency conditions (4.10) yields

$$R_{\ell}^{\circ} = \tilde{R}_{\ell}^{\circ}, \quad S_{1,\ell}^{\circ} = \tilde{S}_{1,\ell}^{\circ}, \quad S_{2,\ell-1}^{\circ} = \tilde{S}_{2,\ell-1}^{\circ}, \quad \ell = 2, \dots, N.$$

As the scheme (4.1) satisfies Dahlquist's root condition, it is numerically stable in the mean-square sense. Now the assertion follows from the estimates (4.9) by means of the stability inequality.

Proof. (of Lemma 4.1) To derive a representation of the local error in the form (4.8) we evaluate and resume the deterministic parts at the point $(X(t_{\ell-2}), t_{\ell-2})$ and separate the stochastic terms carefully over the different subintervals $[t_{\ell-2}, t_{\ell-1}]$ and $[t_{\ell-1}, t_{\ell}]$. This ensures the independence of the random variables. It does make the calculations more messy, though. By rewriting

$$\sum_{j=0}^{2} \alpha_{j} X(t_{\ell-j}) = \alpha_{0} \big(X(t_{\ell}) - X(t_{\ell-1}) \big) + (\alpha_{0} + \alpha_{1}) \big(X(t_{\ell-1}) - X(t_{\ell-2}) \big) + \big(\sum_{j=0}^{2} \alpha_{j} \big) X(t_{\ell-2}),$$

we can express the local error (2.3) as

$$L_{\ell} = \alpha_0 \left(X(t_{\ell}) - X(t_{\ell-1}) \right) + (\alpha_0 + \alpha_1) \left(X(t_{\ell-1}) - X(t_{\ell-2}) \right) + \sum_{j=0}^2 \alpha_j X(t_{\ell-2})$$
$$-h \sum_{j=0}^2 \beta_j f(X(t_{\ell-j}), t_{\ell-j}) - \sum_{j=1}^2 \gamma_j \sum_{r=1}^m g_r(X(t_{\ell-j}), t_{\ell-j}) I_r^{t_{\ell-j}, t_{\ell-j+1}}.$$

The SDE (1.1) implies the identities

$$X(t_{\ell-1}) - X(t_{\ell-2}) = hf(X(t_{\ell-2}), t_{\ell-2}) + I_{00}^{t_{\ell-2}, t_{\ell-1}}(\Lambda_0 f) + \sum_{r=1}^m I_{r0}^{t_{\ell-2}, t_{\ell-1}}(\Lambda_r f) + \sum_{r=1}^m g_r(X(t_{\ell-2}), t_{\ell-2})I_r^{t_{\ell-2}, t_{\ell-1}} + \sum_{r=1}^m I_{0r}^{t_{\ell-2}, t_{\ell-1}}(\Lambda_0 g_r) + \sum_{r,q=1}^m I_{qr}^{t_{\ell-2}, t_{\ell-1}}(\Lambda_q g_r),$$

and, additionally using (4.6),

$$\begin{aligned} X(t_{\ell}) - X(t_{\ell-1}) &= h \Big\{ f(X(t_{\ell-2}), t_{\ell-2}) + I_0^{t_{\ell-2}, t_{\ell-1}}(\Lambda_0 f) + \sum_{r=1}^m I_r^{t_{\ell-2}, t_{\ell-1}}(\Lambda_r f) \Big\} \\ &+ I_{00}^{t_{\ell-1}, t_{\ell}}(\Lambda_0 f) + \sum_{r=1}^m I_{r0}^{t_{\ell-1}, t_{\ell}}(\Lambda_r f) \\ &+ \sum_{r=1}^m g_r(X(t_{\ell-1}), t_{\ell-1}) I_r^{t_{\ell-1}, t_{\ell}} + \sum_{r=1}^m I_{0r}^{t_{\ell-1}, t_{\ell}}(\Lambda_0 g_r) + \sum_{r, q=1}^m I_{qr}^{t_{\ell-1}, t_{\ell}}(\Lambda_q g_r). \end{aligned}$$

Inserting this and the expansions (4.6, 4.7) into the local error formula and reordering the terms, yields

$$L_{\ell} = \left[\sum_{j=0}^{2} \alpha_{j}\right] X(t_{\ell-2}) + \left[2\alpha_{0} + \alpha_{1} - \sum_{j=0}^{2} \beta_{j}\right] hf(X(t_{\ell-2}), t_{\ell-2}) + \tilde{R}_{\ell}^{\circ} + \left[\alpha_{0} - \gamma_{1}\right] \sum_{r=1}^{m} g_{r}(X(t_{\ell-1}), t_{\ell-1}) I_{r}^{t_{\ell-1}, t_{\ell}} + \tilde{S}_{1,\ell}^{\circ} + \left[(\alpha_{0} + \alpha_{1}) - \gamma_{2}\right] \sum_{r=1}^{m} g_{r}(X(t_{\ell-2}), t_{\ell-2}) I_{r}^{t_{\ell-2}, t_{\ell-1}} + \tilde{S}_{2,\ell-1}^{\circ},$$

where

$$\tilde{R}_{\ell}^{\circ} = \alpha_0 \left\{ h I_0^{t_{\ell-2}, t_{\ell-1}}(\Lambda_0 f) + I_{00}^{t_{\ell-1}, t_{\ell}}(\Lambda_0 f) \right\} + (\alpha_0 + \alpha_1) I_{00}^{t_{\ell-2}, t_{\ell-1}}(\Lambda_0 f) - h \beta_0 \left\{ I_0^{t_{\ell-2}, t_{\ell-1}}(\Lambda_0 f) + I_0^{t_{\ell-1}, t_{\ell}}(\Lambda_0 f) \right\} - h \beta_1 I_0^{t_{\ell-2}, t_{\ell-1}}(\Lambda_0 f),$$
(4.11)

$$\tilde{S}_{1,\ell}^{\circ} = \sum_{r=1}^{m} \left(\alpha_0 I_{r0}^{t_{\ell-1},t_{\ell}}(\Lambda_r f) - h\beta_0 I_r^{t_{\ell-1},t_{\ell}}(\Lambda_r f) \right) + \alpha_0 \sum_{r=1}^{m} I_{0r}^{t_{\ell-1},t_{\ell}}(\Lambda_0 g_r) + \alpha_0 \sum_{r,q=1}^{m} I_{qr}^{t_{\ell-1},t_{\ell}}(\Lambda_q g_r)),$$
(4.12)

$$\tilde{S}_{2,\ell-1}^{\circ} = h(\alpha_0 - \beta_0 - \beta_1) \sum_{r=1}^{m} I_r^{t_{\ell-2},t_{\ell-1}}(\Lambda_r f) + (\alpha_0 + \alpha_1) \sum_{r=1}^{m} I_{r0}^{t_{\ell-2},t_{\ell-1}}(\Lambda_r f) + (\alpha_0 + \alpha_1) \sum_{r=1}^{m} I_{0r}^{t_{\ell-2},t_{\ell-1}}(\Lambda_0 g_r) + (\alpha_0 + \alpha_1) \sum_{r,q=1}^{m} I_{qr}^{t_{\ell-2},t_{\ell-1}}(\Lambda_q g_r).$$
(4.13)

Finally, the estimates (4.9) are derived by means of Lemma 2.12, where the last terms in (4.12) and (4.13) determine the order $\mathcal{O}(h)$.

4.2. Two-step schemes for small noise SDEs. For the numerical integration of ODEs two-step schemes of order 2 or higher are particularly interesting. They offer a high order of convergence for low computational cost per step. In this section we discuss the special case of small noise SDEs (1.2) where the parameter ϵ in the diffusion coefficients $g_r = \epsilon \hat{g}_r$, $r = 1, \dots, m$ measures the smallness of the noise. Lemma 4.1 provides a representation for the local error (2.3) of the stochastic linear two step scheme (4.1) applied to (1.2). Starting from this expression we will further analyse the local error by expanding the term $\Lambda_0 f$ appearing in \tilde{R}°_{ℓ} (4.11). Naturally, this requires more smoothness of the coefficients. A sufficient condition would be $\Lambda_0 f \in C^{2,1}$, for which, in general, one needs the existence of fourth order derivatives of f with respect to x. However, for small noise SDEs, the term $f'_x f + f'_t$ dominates $\Lambda_0 f$. This allows weakening the smoothness assumptions again. The expansion of $\Lambda_0 f$ also yields additional multiple Itô-integrals whose conditional expectation vanishes. By moving these terms from \hat{R}°_{ℓ} into the stochastic parts of the representation of the local error we achieve better estimates. With this analysis we are able to prove that some of the potential of deterministic two-step schemes can be recovered in the special case of small noise SDEs.

To be able to exploit the effect of the small parameter ϵ in the expansions of the local error we introduce operators Λ_0^f , $\hat{\Lambda}_0$ and $\hat{\Lambda}_r$, $r = 1, \ldots, m$ defined on $C^{2,1}$ and $C^{1,0}$, respectively, by

$$\Lambda_0^f y := y'_t + y'_x f, \quad \hat{\Lambda}_0 y := \frac{1}{2} \sum_{r=1}^m y''_{xx} [\hat{g}_r, \hat{g}_r], \quad \hat{\Lambda}_r y := y'_x \hat{g}_r .$$
(4.14)

In terms of the original definition (4.2) we have

$$\Lambda_0 y = \Lambda_0^f y + \epsilon^2 \hat{\Lambda}_0 y \quad \text{and} \quad \Lambda_r y = \epsilon \hat{\Lambda}_r y.$$
(4.15)

LEMMA 4.3. Assume that the coefficients $f, \hat{g}_r, r = 1, \ldots, m$ of the small noise SDE (1.2), as well as $\Lambda_0^f f = f'_x f + f'_t$ belong to the class $C^{2,1}$ with $\Lambda_0 f, \Lambda_0 \hat{g}_r, \hat{\Lambda}_r f, \hat{\Lambda}_q \hat{g}_r, \Lambda_0 \Lambda_0^f f, \hat{\Lambda}_r \Lambda_0^f f \in C^K$ for $r, q = 1, \ldots, m$. Let the stochastic 2-step scheme (4.1) satisfy the consistency conditions (4.10). Then the local error (2.3) of the method (4.1) for the small noise SDE (1.2) allows the representation

$$L_{\ell} = R_{\ell}^{\diamond} + S_{1,\ell}^{\diamond} + S_{2,\ell-1}^{\diamond}, \quad \ell = 2, \dots, N,$$
(4.16)

where $R_{\ell}^{\diamond}, S_{j,\ell}^{\diamond}, j = 1, 2$ are $\mathcal{F}_{t_{\ell}}$ -measurable with $\mathbb{E}(S_{j,\ell}^{\diamond}|\mathcal{F}_{t_{\ell-1}}) = 0$, and

$$\begin{aligned} R_{\ell}^{\diamond} &= \left[(4\alpha_0 + \alpha_1) - (4\beta_0 + 2\beta_1) \right] \frac{h^2}{2} (f'_t + f'_x f) (X(t_{\ell-2}), t_{\ell-2}) + \tilde{R}_{\ell}^{\diamond}, \\ S_{1,\ell}^{\diamond} &= \tilde{S}_{1,\ell}^{\circ} + \tilde{S}_{1,\ell}^{\diamond}, \\ S_{2,\ell-1}^{\diamond} &= \tilde{S}_{2,\ell-1}^{\circ} + \tilde{S}_{2,\ell-1}^{\diamond}, \end{aligned}$$

where

$$\|\tilde{R}_{\ell}^{\diamond}\|_{L_{2}} = O(h^{3} + \epsilon^{2}h^{2}), \quad \|\tilde{S}_{1,\ell}^{\diamond}\|_{L_{2}} = O(\epsilon h^{5/2}), \quad \|\tilde{S}_{2,\ell-1}^{\diamond}\|_{L_{2}} = O(\epsilon h^{5/2}).$$
(4.17)

The terms $\tilde{S}_{1,\ell}^{\circ}, \tilde{S}_{2,\ell-1}^{\circ}$ are given by (4.12, 4.13) in the proof of Lemma 4.1 and satisfy here

$$\|\tilde{S}_{1,\ell}^{\circ}\|_{L_2} = \mathcal{O}(\epsilon^2 h + \epsilon h^{3/2}), \quad \|\tilde{S}_{2,\ell}^{\circ}\|_{L_2} = \mathcal{O}(\epsilon^2 h + \epsilon h^{3/2}).$$
(4.18)

Proof. We have from Lemma 4.1, if the consistency conditions (4.10) are satisfied, the representation

$$L_{\ell} = \tilde{R}_{\ell}^{\circ} + \tilde{S}_{1,\ell}^{\circ} + \tilde{S}_{2,\ell-1}^{\circ}, \ \ell = 2, \dots, N_{\ell}$$

where \tilde{R}°_{ℓ} , $\tilde{S}^{\circ}_{1,\ell}$, $\tilde{S}^{\circ}_{2,\ell-1}$ are given by (4.11, 4.12, 4.13). Splitting $\Lambda_0 f = \Lambda_0^f f + \epsilon^2 \hat{\Lambda}_0 f$ immediately yields $\tilde{R}^{\circ}_{\ell} = \tilde{R}^{\circ f}_{\ell} + \epsilon^2 \hat{R}^{\circ}_{\ell}$ with

$$\tilde{R}_{\ell}^{\circ f} := (\alpha_0 - \beta_0 - \beta_1) h I_0^{t_{\ell-2}, t_{\ell-1}} (\Lambda_0^f f) + (\alpha_0 + \alpha_1) I_{00}^{t_{\ell-2}, t_{\ell-1}} (\Lambda_0^f f)
+ \alpha_0 I_{00}^{t_{\ell-1}, t_{\ell}} (\Lambda_0^f f) - h \beta_0 I_0^{t_{\ell-1}, t_{\ell}} (\Lambda_0^f f)$$
(4.19)

$$\hat{R}^{\circ}_{\ell} := (\alpha_0 - \beta_0 - \beta_1) h I_0^{t_{\ell-2}, t_{\ell-1}} (\hat{\Lambda}_0 f) + (\alpha_0 + \alpha_1) I_{00}^{t_{\ell-2}, t_{\ell-1}} (\hat{\Lambda}_0 f) + \alpha_0 I_{00}^{t_{\ell-1}, t_{\ell}} (\hat{\Lambda}_0 f) - h \beta_0 I_0^{t_{\ell-1}, t_{\ell}} (\hat{\Lambda}_0 f).$$
(4.20)

We note that (4.20) appears with the factor ϵ^2 in the local error representation, thus yielding the $\mathcal{O}(\epsilon^2 h^2)$ term in the estimate of $\|\tilde{R}_\ell^{\diamond}\|_{L_2}$ in (4.17). We concentrate on developing $\tilde{R}_\ell^{\circ f}$ in more detail. Applying the Itô-formula (4.3) to $\Lambda_0^f f(X(s),s)$ for $s \in [t_{\ell-2}, t_{\ell-1}]$ and integrating yields

$$I_0^{t_{\ell-2},s}(\Lambda_0^f f) = (s - t_{\ell-2})\Lambda_0^f f(X(t_{\ell-2}), t_{\ell-2}) + I_{00}^{t_{\ell-2},s}(\Lambda_0\Lambda_0^f f) + \epsilon \sum_{r=1}^m I_{r0}^{t_{\ell-2},s}(\hat{\Lambda}_r\Lambda_0^f f).$$

For $s=t_{\ell-1}$ we obtain an expression for the first integral in (4.19) . Integrating again we obtain for the second integral in (4.19)

$$I_{00}^{t_{\ell-2},t_{\ell-1}}(\Lambda_0^f f) = \frac{h^2}{2}\Lambda_0^f f(X(t_{\ell-2}),t_{\ell-2}) + I_{000}^{t_{\ell-2},t_{\ell-1}}(\Lambda_0\Lambda_0^f f) + \epsilon \sum_{r=1}^m I_{r00}^{t_{\ell-2},t_{\ell-1}}(\hat{\Lambda}_r\Lambda_0^f f)$$

Both the other integrals are over the interval $[t_{\ell-1}, t_{\ell}]$. In the analogous expressions for these the term $\Lambda_0^f f(X(t_{\ell-1}), t_{\ell-1})$ has to be substituted by

 $\Lambda_0^f f(X(t_{\ell-1}), t_{\ell-1}) = \Lambda_0^f f(X(t_{\ell-2}), t_{\ell-2}) + I_0^{t_{\ell-2}, t_{\ell-1}} (\Lambda_0 \Lambda_0^f f) + \epsilon \sum_{r=1}^m I_r^{t_{\ell-2}, t_{\ell-1}} (\Lambda_r \Lambda_0^f f).$ Then we obtain from (4.19)

$$\tilde{R}_{\ell}^{\circ f} = \left[(4\alpha_0 + \alpha_1) - (4\beta_0 + 2\beta_1) \right] \frac{h^2}{2} \Lambda_0^f f(X(t_{\ell-2}), t_{\ell-2}) + \tilde{R}_{\ell}^{\diamond f} + \tilde{S}_{1,\ell}^{\diamond} + \tilde{S}_{2,\ell}^{\diamond},$$

where

$$\tilde{R}_{\ell}^{\diamond f} = (\alpha_0 - 2\beta_0) \frac{h^2}{2} I_0^{t_{\ell-2}, t_{\ell-1}} (\Lambda_0 \Lambda_0^f f) + (\alpha_0 - \beta_0 - \beta_1) h I_{00}^{t_{\ell-2}, t_{\ell-1}} (\Lambda_0 \Lambda_0^f f) - \beta_0 h I_{00}^{t_{\ell-1}, t_{\ell}} (\Lambda_0 \Lambda_0^f f) + (\alpha_0 + \alpha_1) I_{000}^{t_{\ell-2}, t_{\ell-1}} (\Lambda_0 \Lambda_0^f f) + \alpha_0 I_{000}^{t_{\ell-1}, t_{\ell}} (\Lambda_0 \Lambda_0^f f),$$

$$\begin{split} \tilde{S}_{1,\ell}^{\diamond} &= \alpha_0 \epsilon \sum_{r=1}^m I_{r00}^{t_{\ell-1},t_{\ell}} (\hat{\Lambda}_r \Lambda_0^f f) - h\beta_0 \epsilon \sum_{r=1}^m I_{r0}^{t_{\ell-1},t_{\ell}} (\hat{\Lambda}_r \Lambda_0^f f), \\ \tilde{S}_{2,\ell}^{\diamond} &= (\alpha_0 - 2\beta_0) \frac{h^2}{2} \epsilon \sum_{r=1}^m I_r^{t_{\ell-2},t_{\ell-1}} (\hat{\Lambda}_r \Lambda_0^f f) + (\alpha_0 - \beta_0 - \beta_1) h\epsilon \sum_{r=1}^m I_{r0}^{t_{\ell-2},t_{\ell-1}} (\hat{\Lambda}_r \Lambda_0^f f) \\ &+ (\alpha_0 + \alpha_1) \epsilon \sum_{r=1}^m I_{r00}^{t_{\ell-2},t_{\ell-1}} (\hat{\Lambda}_r \Lambda_0^f f). \end{split}$$

We arrive at $\tilde{R}_{\ell}^{\diamond} = \tilde{R}_{\ell}^{\diamond f} + \epsilon^2 \hat{R}_{\ell}^{\diamond}$. Finally, the estimates (4.17) are derived by means of Lemma 2.12.

COROLLARY 4.4. Let the coefficients f, \hat{g}_r , r = 1, ..., m, of the SDE (1.2) satisfy the assumptions of Lemma 4.3 and suppose they are Lipschitz continuous with respect to their first variable. Let the coefficients of the stochastic linear two-step scheme (4.1) satisfy Dahlquist's root condition and the consistency conditions (4.10) and

$$(4\alpha_0 + \alpha_1) - (4\beta_0 + 2\beta_1) = 0. \tag{4.21}$$

Then the global error of the scheme (4.1) applied to (1.2) allows the expansion

$$\max_{\ell=2,\dots,N} \|X(t_{\ell}) - X_{\ell}\|_{L_{2}} = \mathcal{O}(h^{2} + \epsilon h + \epsilon^{2} h^{1/2}) + \mathcal{O}(\max_{\ell=0,1} \|X(t_{\ell}) - X_{\ell}\|_{L_{2}}).$$

Proof. Lemma 4.3 stated the representation (4.16) for the local error. Applying the consistency condition (4.21) yields $R_{\ell}^{\diamond} = \tilde{R}_{\ell}^{\diamond}$ and by (4.17) we have $||R_{\ell}^{\diamond}||_{L_2} = \mathcal{O}(h^3 + \epsilon^2 h^2)$. The stochastic terms $S_{1,\ell}^{\diamond}$, $S_{2,\ell-1}^{\diamond}$ are dominated by $\tilde{S}_{1,\ell}^{\circ}$, $\tilde{S}_{2,\ell-1}^{\circ}$ and thus are of order of magnitude $\mathcal{O}(\epsilon^2 h + \epsilon h^{3/2})$. As the scheme (4.1) satisfies Dahlquist's root condition, it is numerically stable in the mean-square sense. Applying the stability inequality (2.8) to the representation (4.16) of the local error yields the assertion.

REMARK 4.5. We note that R_{ℓ}^{\diamond} is responsible for the $\mathcal{O}(h^2)$ term in the expansion of the global error. In the limit $\epsilon \to 0$ this is the only remaining term. It reflects the convergence properties of the scheme in the deterministic setting. On the other hand, for asymptotically small stepsizes $h \to 0$ and fixed parameter ϵ , even if it is small, the term of order $\mathcal{O}(\epsilon^2 h^{1/2})$ causes the low order $\frac{1}{2}$ of convergence. The question arises for which choices of ϵ and h the schemes still show the order 2 behaviour. Thus we are interested in when the term $\mathcal{O}(h^2)$ dominates the term $\mathcal{O}(\epsilon h + \epsilon^2 h^{1/2})$. Clearly both these terms depend on the actual coefficients $f, \hat{g}_r, r = 1, \ldots, m$, of the SDE (1.2) and their derivatives. Assuming moderate function values, the term $\mathcal{O}(h^2)$ dominates $\mathcal{O}(\epsilon h + \epsilon^2 h^{1/2})$, if $h^2 \gg \epsilon^2 h^{1/2}$, i.e. $h \gg \epsilon^{4/3}$, and $h^2 \gg \epsilon h$, i.e. $h \gg \epsilon$. Obviously, in general the second condition is stronger. Summarizing, we can expect order 2 behaviour if $h \gg \epsilon$.

However, even if the chosen stepsize does not satisfy this condition for a given ϵ , schemes satisfying the consistency condition (4.21) often show a better performance than other schemes. The reason is that their error is dominated by $\mathcal{O}(\epsilon h)$ instead of $\mathcal{O}(\epsilon^2 h^{1/2})$, resulting in an order 1 behaviour with the small parameter ϵ in the error constant. Again assuming moderate function values, one may expect this for $\epsilon h \gg \epsilon^2 h^{1/2}$, i.e. $h \gg \epsilon^2$.

5. Test results. We report results for several explicit and implicit stochastic linear k-step schemes for k = 1, 2, applied to two examples of SDEs. Table 5.1 summarizes the methods we have implemented and tested. All methods, including the one-step schemes, satisfy the consistency conditions (4.10) and all methods, excluding the Euler schemes, satisfy the consistency condition (4.21). Further, all methods, exept the first one, are zero-stable. The Milne-Simpson method is the only linear two-step method with deterministic order of convergence 4. However, its characteristic polynomial possesses the second root $\alpha_2 = -1$ on the unit circle which easily causes weak instabilities in the integration of decaying solutions, such that this scheme is generally not recommended.

Method	α_0	α_1	α_2	eta_0	β_1	β_2	γ_1	γ_2
unstable method	1	-3	2	0	1/2	-3/2	1	-2
explicit Euler	1	-1	0	0	1	0	1	0
implicit Euler	1	-1	0	1	0	0	1	0
trapezoidal rule	1	-1	0	1/2	1/2	0	1	0
BDF 2	1	-4/3	1/3	2/3	0	0	1	-1/3
Adams-Bashforth 2	1	-1	0	0	3/2	-1/2	1	0
Adams-Moulton 2	1	-1	0	5/12	8/12	-1/12	1	0
Milne-Simpson	1	0	-1	1/3	4/3	1/3	1	1

TABLE 5.1

Coefficients of two-step schemes

The first example is the simple bilinear scalar SDE

$$X(t) = 1 + \int_0^t \alpha X(s) \, \mathrm{d}s + \int_0^t \beta X(s) \, \mathrm{d}W(s) \, , \, t \in [0, 1],$$
 (5.1)

with coefficients $f(x,t) := \alpha x$, $G(x,t) = (g(x,t)) = (\beta x)$, parameters $\alpha, \beta \in \mathbb{R}$ and a scalar Wiener process W. This example has the advantage that an explicit solution is available to compare with the numerical approximations. This solution is given by the geometric Brownian motion $X(t) = \exp\left((\alpha - \frac{1}{2}\beta^2)t + \beta W(t)\right)$.

The second example is a two-dimensional system of SDEs which is taken from [8, 9, Problem P1], where it is given in Stratonovich notation. For the purpose of using it as a test equation for small noise SDEs in the form (1.2) we have transformed it into Itô notation and included a parameter ϵ . We arrive at

$$X(t) = X(0) + \int_0^t FX(s) \, ds + \int_0^t G_1 X(s) \, dW_1(s) + \int_0^t G_2 X(s) \, dW_2(s), \quad (5.2)$$

$$F = \begin{pmatrix} -0.9 & 0 \\ 0.25 & -0.5 \end{pmatrix} + \frac{\epsilon^2}{2} \begin{pmatrix} 0.75^2 + 0.9^2 & 0 \\ 0 & 0.75^2 + 0.9^2 \end{pmatrix},$$

$$G_1 = \epsilon \begin{pmatrix} 0.75 & 0 \\ 0 & -0.75 \end{pmatrix}, \quad G_2 = \epsilon \begin{pmatrix} 0 & 0.9 \\ 0.9 & 0 \end{pmatrix}, \quad t \in [0, 2].$$

The above system is fully non-commutative, i.e., $[F, G_1]$, $[F, G_2]$, $[G_1, G_2]$ are all nonzero, where [A, B] = AB - BA for matrices A, B. Here we can not use an explicit formula for the solution of the system, thus we have computed a 'reference solution' with the trapezoidal rule on a very fine grid by using 262144 steps.

Results for another example, a scalar SDE with polynomial drift and diffusion coefficients, are reported in [6].

In our experiments we have investigated the relation between step-size h and achieved accuracy for several choices of parameters for the two test examples. The accuracy is measured as the maximum approximate L_2 -norm of the global errors in the timeinterval [0, 1] and [0, 2], respectively:

$$\max_{\ell=1,\dots,N} \left(\frac{1}{M} \sum_{j=1}^{M} |X(t_{\ell},\omega_j) - X_{\ell}(\omega_j)|^2 \right)^{1/2} \approx \max_{\ell=1,\dots,N} \|X(t_{\ell}) - X_{\ell}\|_{L_2}$$

where N denotes the number of steps and M the number of computed paths. In our computations we used M = 100, unless specified otherwise.

The results are presented as figures, where we have plotted the achieved accuracy versus the step-sizes in logarithmic scale with base 10. Then the slope of the resulting lines corresponds to the observed order of the schemes. Lines with slopes 0.5, 1, 2, 3 are provided in some figures to enable comparisons with convergence of these orders. Our first test concerned the effect of applying a method which is not numerically stable in the mean-square sense, i.e. the coefficients of the method do not satisfy Dahlquist's root condition. For the method considered (the parameters are given in the first row of 5.1), one of the roots of the characteristic polynomial ρ (2.12) is 2. Figure 1 shows the behaviour of the error of this method applied to the SDE (5.1) with $\alpha = -1$, $\beta = 0.01$, when the step-size decreases. We note that the scale is logarithmic!

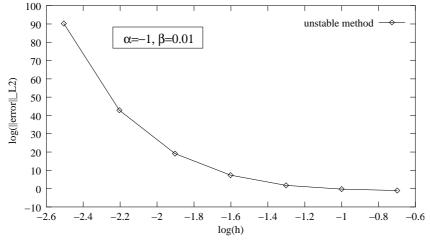


Figure 1: Simulation results for the SDE (5.1) with the unstable scheme

Next we report results for the other methods listed in Table 5.1, applied to the simple linear SDE (5.1). To start off the integration with the two-step schemes we needed a second starting value X_1 . In this example we used the exact solution value $X(t_1)$, thus avoiding introducing additional errors. In computational practice, the starting value X_1 could be computed, e.g., by means of the trapezoidal rule. Then the error structure of the two-step Maruyama schemes is maintained. For reasons of space restriction we present only results for one set of parameters, $\alpha = -1$ and $\beta = 0.01$, where β takes the role of the parameter ϵ . The results are given in Figure 2 and illustrate the observations made in Remark 4.5 very well. We see that the error of the other schemes is smaller than that of the Euler schemes. Their error appears to be of the size $\max(c_1\beta^2h^{\frac{1}{2}}, c_2\beta h, c_3h^2)$, c_3 being an error constant particular to the scheme. The constants c_1, c_2 appear to coincide for the considered schemes.

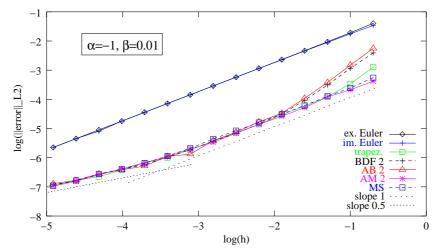


Figure 2: Simulation results for the SDE (5.1) with the stable schemes of Table 5.1

Results of experiments with several methods of Table 5.1 applied to the second example (5.2) are presented in Figures 3 to 5. The second starting value X_1 has been taken as the corresponding value of the reference solution, which amounts to a computation by the trapezoidal rule with very small step-sizes. In Figures 3 and 4 we show simulation results with the parameter $\epsilon = 10^{-4}$, 10^{-2} . Again we have plotted the achieved accuracy versus the step-size in logarithmic scale. The results are qualitatively the same as for the simple example above: The error of the Euler schemes is larger than that of the other schemes and the error of the schemes with deterministic order 2 (trapezoidal rule, Adams-Bashforth2 and BDF2) appears to be of the size $\max(c_1\beta^2h^{\frac{1}{2}}, c_2\beta h, c_3h^2)$, whereas for the Adams-Moulton2 with deterministic order 3 one even observes errors of the size $\max(c_1\beta^2h^{\frac{1}{2}}, c_2\beta h, c_3h^2)$.

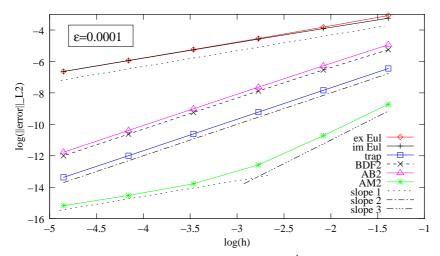


Figure 3: Simulation results for the SDE (5.2), $\epsilon = 10^{-4}$, with schemes of Table 5.1

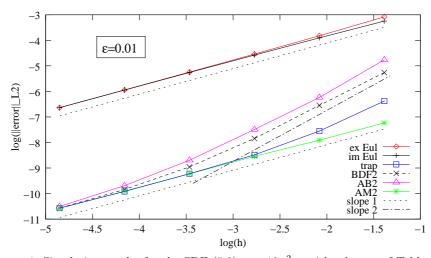


Figure 4: Simulation results for the SDE (5.2), $\epsilon = 10^{-2}$, with schemes of Table 5.1 Figure 5 relates to the case where the parameter in the diffusion coefficient is not small. Here the order $h^{1/2}$ term dominates the error for all chosen step-sizes. Furthermore, due to the larger noise the statistical error in the approximation is more visible, which we partly compensated by using 500 paths instead of only 100. In this case one can clearly see that the Euler schemes perform as good as the other schemes and investing the higher effort for the two-step schemes does not pay.

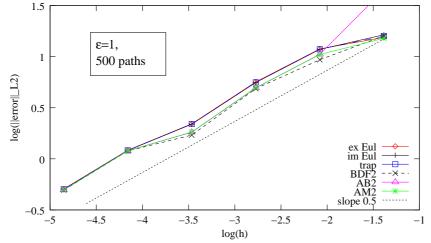


Figure 5: Simulation results for the SDE (5.2), $\epsilon = 1$, with schemes of Table 5.1

Appendix A. Proof of Theorem 3.2.

Proof. Necessity: This part can be proved as in the deterministic case, i. e. we take the equation X'(t) = 0, then f and Γ_j satisfy obviously (2.9) and (2.10). We then follow in principle the proof of [15, Thm.6.3.3].

Sufficiency: Since the SLMM (2.1) contains the stochastic part related to the Γ_j , we can not rely on the theory of difference equations and the representations of their solutions. Instead, we will follow the route of rewriting the *k*-step recurrence equation as a one-step recurrence equation in a higher dimensional space (see e.g. [16, Chap.III.4][26, Chap.8.2.1]).

For X_{ℓ} and \tilde{X}_{ℓ} being the solutions of (2.1) and (2.4), respectively, let the *n*-dimensional vector E_{ℓ} be defined as the difference $X_{\ell} - \tilde{X}_{\ell}$. We have with $E_0, \ldots, E_{k-1} \in L_2(\Omega, \mathbb{R}^n)$

for $\ell = k, \ldots, N$, the recursion

$$E_{\ell} = -\sum_{j=1}^{k} \alpha_{j} \ E_{\ell-j} + \underbrace{h \ \sum_{j=0}^{k} \beta_{j} \ \Delta f_{\ell-j}}_{=:\Delta \phi^{\ell}} + \underbrace{\sum_{j=1}^{k} \Delta \Gamma_{j,\ell-j} \ I^{t_{\ell-j},t_{\ell-j+1}}}_{=:\Delta \psi^{\ell}} - D_{\ell}, \qquad (A.1)$$

where

$$\Delta f_{\ell-j} := f(X_{\ell-j}, t_{\ell-j}) - f(X_{\ell-j}, t_{\ell-j})$$

$$\Delta \Gamma_{j,\ell-j} := \Gamma_j(X_{\ell-j}, t_{\ell-j}) - \Gamma_j(\tilde{X}_{\ell-j}, t_{\ell-j}).$$

We rearrange this k-step recursion in the space $L_2(\Omega, \mathbb{R}^n)$ to a one-step recursion in $L_2(\Omega, \mathbb{R}^{k \times n})$. Together with the trivial identities $E_{\ell-1} = E_{\ell-1}, \ldots E_{\ell-k+1} = E_{\ell-k+1}$ we obtain

$$\underbrace{\begin{pmatrix} E_{\ell} \\ E_{\ell-1} \\ \vdots \\ E_{\ell-k+1} \end{pmatrix}}_{=: \mathcal{E}_{\ell}} = \underbrace{\begin{pmatrix} -\alpha_1 I \cdots - \alpha_k I \\ I & 0 \\ & \ddots & \ddots \\ & I & 0 \end{pmatrix}}_{=: \mathcal{A}} \underbrace{\begin{pmatrix} E_{\ell-1} \\ E_{\ell-2} \\ \vdots \\ E_{\ell-k} \end{pmatrix}}_{=: \mathcal{E}_{\ell-1}} + \underbrace{\begin{pmatrix} \Delta \phi^{\ell} \\ 0 \\ \vdots \\ 0 \end{pmatrix}}_{=: \Delta \Phi_{\ell}} + \underbrace{\begin{pmatrix} \Delta \psi^{\ell} \\ 0 \\ \vdots \\ 0 \end{pmatrix}}_{=: \Delta \Psi_{\ell}} + \underbrace{\begin{pmatrix} -D_{\ell} \\ 0 \\ \vdots \\ 0 \end{pmatrix}}_{=: \Delta \Psi_{\ell}} + \underbrace{\begin{pmatrix} -D_{\ell} \\ 0 \\ \vdots \\ 0 \end{pmatrix}}_{=: \mathcal{D}_{\ell}}$$

or, in compact form

$$\mathcal{E}_{\ell} = \mathcal{A}\mathcal{E}_{\ell-1} + \Delta \Phi_{\ell} + \Delta \Psi_{\ell} + \mathcal{D}_{\ell} , \ \ell = k, \dots, N \quad \text{and} \quad \mathcal{E}_{k-1} = (-D_{k-1}, -D_{k-2}, \dots, -D_0)^T,$$

where $\mathcal{E}_{\ell} \in L_2(\Omega, \mathbb{R}^{k \times n})$, $\ell = k - 1, k, \ldots, N$. The vector \mathcal{E}_{k-1} consists of the perturbations to the initial values. We now trace back the recursion in \mathcal{E}_{ℓ} to the initial vector \mathcal{E}_{k-1} . For $\ell = k, \ldots, N$ we have

$$\begin{split} \mathcal{E}_{\ell} &= \mathcal{A}\mathcal{E}_{\ell-1} + \Delta\Phi_{\ell} + \Delta\Psi_{\ell} + \mathcal{D}_{\ell} \\ &= \mathcal{A}(\mathcal{A}\mathcal{E}_{\ell-2} + \Delta\Phi_{\ell-1} + \Delta\Psi_{\ell-1} + \mathcal{D}_{\ell-1}) + \Delta\Phi_{\ell} + \Delta\Psi_{\ell} + \mathcal{D}_{\ell} \\ &= \mathcal{A}^{2}\mathcal{E}_{\ell-2} + (\Delta\Phi_{\ell} + \mathcal{A}\Delta\Phi_{\ell-1}) + (\Delta\Psi_{\ell} + \mathcal{A}\Delta\Psi_{\ell-1}) + (\mathcal{D}_{\ell} + \mathcal{A}\mathcal{D}_{\ell-1}) \\ &\vdots \\ &= \mathcal{A}^{\ell-k+1}\mathcal{E}_{k-1} + \sum_{i=0}^{\ell-k}\mathcal{A}^{i}\Delta\Phi_{\ell-i} + \sum_{i=0}^{\ell-k}\mathcal{A}^{i}\Delta\Psi_{\ell-i} + \sum_{i=0}^{\ell-k}\mathcal{A}^{i}\mathcal{D}_{\ell-i} \\ &= \mathcal{A}^{\ell-k+1}\mathcal{E}_{k-1} + \sum_{i=k}^{\ell}\mathcal{A}^{\ell-i}\Delta\Phi_{i} + \sum_{i=k}^{\ell}\mathcal{A}^{\ell-i}\Delta\Psi_{i} + \sum_{i=k}^{\ell}\mathcal{A}^{\ell-i}\mathcal{D}_{i} \;. \end{split}$$

A crucial point for the subsequent calculations is to find a scalar product inducing a matrix norm such that this norm of the matrix \mathcal{A} is less than or equal to 1 (see e.g. [16, Chap.III.4,Lemma 4.4]. This is possible if the eigenvalues of the Frobenius matrix \mathcal{A} lie inside the unit circle of the complex plane and are simple if their modulus is equal to 1. The eigenvalues of \mathcal{A} are the roots of the characteristic polynomial ρ (2.12) and due to the assumption that Dahlquist's root condition is satisfied they have the required property. Then there exists a non-singular matrix \mathcal{C} with a block-structure like \mathcal{A} such that $\|\mathcal{C}^{-1}\mathcal{A}\mathcal{C}\|_2 \leq 1$, where $\|\cdot\|_2$ denotes the spectral matrix norm that is induced by the Euclidian vector norm in $\mathbb{R}^{k \times n}$. We can thus choose a scalar product for $\mathcal{X}, \mathcal{Y} \in \mathbb{R}^{k \times n}$ as

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$$\langle \mathcal{X}, \mathcal{Y} \rangle_* := \langle \mathcal{C}^{-1} \mathcal{X}, \mathcal{C}^{-1} \mathcal{Y} \rangle_2$$

and then have $|.|_*$ as the induced vector norm on $\mathbb{R}^{k \times n}$ and $\|\cdot\|_*$ as the induced matrix norm with $\|\mathcal{A}\|_* = \|\mathcal{C}^{-1}\mathcal{A}\mathcal{C}\|_2 \leq 1$. We also have

$$\langle \mathcal{X}, \mathcal{Y} \rangle_* = \mathcal{X}^T \mathcal{C}^{-T} \mathcal{C}^{-1} \mathcal{Y} = \mathcal{X}^T \mathcal{C}^* \mathcal{Y} \text{ with } \mathcal{C}^* = \mathcal{C}^{-T} \mathcal{C}^{-1} = (c_{ij}^* I_n)_{i,j=1,\dots,k}$$

Due to the norm equivalence there are constants $c^{\ast},c_{\ast}>0$ such that

$$|\mathcal{X}|_2^2 \le c^* |\mathcal{X}|_*^2$$
 and $|\mathcal{X}|_*^2 \le c_* |\mathcal{X}|_\infty^2$ $\forall \mathcal{X} \in \mathbb{R}^{k \times n}$

where $|\mathcal{X}|_2^2 = \sum_{j=1,\dots,k} |x_j|^2$, $|\mathcal{X}|_{\infty} = \max_{j=1,\dots,k} |x_j|$ for $\mathcal{X} = (x_1^T,\dots,x_k^T)^T$. For the special vectors $\mathcal{X} = (x^T, 0, \dots, 0)^T$ and $\mathcal{Y} = (y^T, 0, \dots, 0)^T$ with $\mathcal{X}, \mathcal{Y} \in \mathbb{R}^{k \times n}$ and $x, y \in \mathbb{R}^n$, one has $\langle \mathcal{X}, \mathcal{Y} \rangle_* = c_{11}^* \langle x, y \rangle_2 = c_{11}^* x^T y$, where c_{11}^* is given by the matrix \mathcal{C}^* .

We now apply $|.|_*^2$ to estimate $|\mathcal{E}_\ell|_*^2$ and, later, $\mathbb{E}|\mathcal{E}_\ell|_*^2$. We start with

$$|\mathcal{E}_{\ell}|_{*}^{2} \leq 4 \bigg\{ \underbrace{|\mathcal{A}^{\ell-k+1}\mathcal{E}_{k-1}|_{*}^{2}}_{1)} + \underbrace{|\sum_{i=k}^{\ell} \mathcal{A}^{\ell-i} \Delta \Phi_{i}|_{*}^{2}}_{2)} + \underbrace{|\sum_{i=k}^{\ell} \mathcal{A}^{\ell-i} \Delta \Psi_{i}|_{*}^{2}}_{3)} + \underbrace{|\sum_{i=k}^{\ell} \mathcal{A}^{\ell-i} \mathcal{D}_{i}|_{*}^{2}}_{4)} \bigg\}.$$

For the term labelled 1) we have $|\mathcal{A}^{\ell-k+1}\mathcal{E}_{k-1}|_*^2 \leq |\mathcal{E}_{k-1}|_*^2$, and thus

$$\mathbb{E}|\mathcal{A}^{\ell-k+1}\mathcal{E}_{k-1}|_*^2 \le \mathbb{E}|\mathcal{E}_{k-1}|_*^2 .$$
(A.2)

For the term labelled 2) we have

$$\begin{split} &|\sum_{i=k}^{\ell} \mathcal{A}^{\ell-i} \Delta \Phi_i|_*^2 \leq (\ell-k+1) \sum_{i=k}^{\ell} |\mathcal{A}^{\ell-i} \Delta \Phi_i|_*^2 \leq N \sum_{i=k}^{\ell} |\Delta \Phi_i|_*^2 = \frac{T}{h} c_{11}^* \sum_{i=k}^{\ell} |\Delta \phi^i|^2 \\ &\leq h \ T \ c_{11}^* \ (k+1) \sum_{i=k}^{\ell} \sum_{j=0}^{k} |\beta_j \ \Delta f_{i-j}|^2 \leq h \ T \ c_{11}^* \ (k+1) \ L_f^2 \sum_{i=k}^{\ell} \sum_{j=0}^{k} \beta_j^2 \ |E_{i-j}|^2 \\ &\leq h \ T \ c_{11}^* \ (k+1) \ L_f^2 \Big\{ \beta_0^2 \ |E_\ell|^2 + \sum_{i=k}^{\ell} \{\beta_0^2 \ |E_{i-1}|^2 + \sum_{j=1}^{k} \beta_j^2 \ |E_{i-j}|^2 \Big\} \Big\} \\ &\leq h \ T \ c_{11}^* \ (k+1) \ L_f^2 \Big\{ c^* \ \beta_0^2 \ |\mathcal{E}_\ell|_*^2 + C_\beta \ c^* \ \sum_{i=k-1}^{\ell-1} |\mathcal{E}_i|_*^2 \Big\}, \end{split}$$

where $C_{\beta} = 2 \max_{j=0,...,k} \beta_j$. Hence,

$$\mathbb{E}|\sum_{i=k}^{\ell} \mathcal{A}^{\ell-i} \Delta \Phi_i|_*^2 \le h \ T \ c_{11}^* \ (k+1) \ L_f^2 \left\{ c^* \ \beta_0^2 \ \mathbb{E}|\mathcal{E}_\ell|_*^2 + C_\beta \ c^* \sum_{i=k-1}^{\ell-1} \mathbb{E}|\mathcal{E}_i|_*^2 \right\}.$$
(A.3)

We will now treat the term labelled 3). For that purpose we introduce the notation $\Delta \Psi_{j,i-j} := ((\Delta \Gamma_{j,i-j} I^{t_{i-j},t_{i-j+1}})^T, 0, \dots, 0)^T$. Using this we can write

$$\Delta \Psi_{i} = ((\Delta \psi^{i})^{T}, 0, \dots, 0)^{T} = ((\sum_{j=1}^{k} \Delta \Gamma_{j,i-j} \ I^{t_{i-j}, t_{i-j+1}})^{T}, 0, \dots, 0)^{T} = \sum_{j=1}^{k} \Delta \Psi_{j,i-j}$$

and $|\sum_{i=k}^{\ell} \mathcal{A}^{\ell-i} \Delta \Psi_{i}|_{*}^{2} = |\sum_{i=k}^{\ell} \mathcal{A}^{\ell-i} \sum_{j=1}^{k} \Delta \Psi_{j,i-j}|_{*}^{2}.$

Every $\Delta \Psi_{j,i-j}$ is $\mathcal{F}_{t_{i-j+1}}$ -measurable and $\mathbb{E}(\Delta \Psi_{j,i-j} | \mathcal{F}_{t_{i-j}}) = 0$. We can now reorder the last term above such that we have a sum of terms where each term contains all multiple Wiener integrals over just one subinterval. The expectation of products of terms from different subintervals vanishes, hence we obtain

$$\begin{split} \mathbb{E} | \sum_{i=k}^{\ell} \mathcal{A}^{\ell-i} \Delta \Psi_{i} |_{*}^{2} \\ &= \mathbb{E} | \mathcal{A}^{\ell-k} \Delta \Psi_{k,0} |_{*}^{2} \\ &+ \mathbb{E} | \mathcal{A}^{\ell-k-1} \Delta \Psi_{k,1} + \mathcal{A}^{\ell-k} \Delta \Psi_{k-1,1} |_{*}^{2} \\ &\vdots \\ &+ \mathbb{E} | \mathcal{A}^{\ell-2k+1} \Delta \Psi_{k,k-1} + \mathcal{A}^{\ell-2k+2} \Delta \Psi_{k-1,k-1} + \dots + \mathcal{A}^{\ell-k} \Delta \Psi_{1,k-1} |_{*}^{2} \\ &\vdots \\ &+ \mathbb{E} | \mathcal{A}^{0} \Delta \Psi_{k,\ell-k} + \mathcal{A}^{1} \Delta \Psi_{k-1,\ell-k} + \dots + \mathcal{A}^{k-1} \Delta \Psi_{1,\ell-k} |_{*}^{2} \\ &\vdots \\ &+ \mathbb{E} | \mathcal{A}^{0} \Delta \Psi_{2,\ell-2} + \mathcal{A}^{1} \Delta \Psi_{1,\ell-2} |_{*}^{2} \\ &+ \mathbb{E} | \mathcal{A}^{0} \Delta \Psi_{1,\ell-1} |_{*}^{2} \\ &\leq k \sum_{i=k}^{\ell} \sum_{j=1}^{k} \mathbb{E} | \Delta \Psi_{j,i-j} |_{*}^{2} \qquad \leq k c_{11}^{*} \sum_{i=k}^{\ell} \sum_{j=1}^{k} \mathbb{E} | | \Delta \Gamma_{j,i-j} ||^{2} \mathbb{E} | I^{t_{i-j},t_{i-j+1}} |^{2} \\ &\leq h \ k \ c_{11}^{*} \ L_{\Gamma}^{2} \sum_{i=k}^{\ell} \sum_{j=1}^{k} \mathbb{E} | \mathbb{E}_{i-j} |^{2} \ \leq h \ k \ c_{11}^{*} \ L_{\Gamma}^{2} \ c^{*} \ \sum_{i=k}^{\ell} | \mathcal{E}_{i-1} |_{*}^{2}. \end{split}$$

Thus, for the term labelled 3), we obtain

$$\mathbb{E}|\sum_{i=k}^{\ell} \mathcal{A}^{\ell-i} \Delta \Psi_i|_*^2 \leq h \ k \ c_{11}^* \ L_{\Gamma}^2 \ c^* \ \sum_{i=k-1}^{\ell-1} |\mathcal{E}_i|_*^2.$$
(A.4)

We will, for a shorter notation, deal with the term labelled 4), i.e. the perturbations D_i in \mathcal{D}_i , after obtaining an intermediate result. Using (A.2), (A.3) and (A.4) and setting $L_0 := L_f^2 (k+1) c_{11}^* T c^* \beta_0^2$ and $L := L_f^2 (k+1) c_{11}^* T c^*_\beta + L_{\Gamma}^2 k c_{11}^* c^*$, we have now arrived at

$$\mathbb{E}|\mathcal{E}_{\ell}|_{*}^{2} \leq 4\left\{\mathbb{E}|\mathcal{E}_{k-1}|_{*}^{2} + hL_{0}\mathbb{E}|\mathcal{E}_{\ell}|^{2} + hL\sum_{i=k-1}^{\ell-1}\mathbb{E}|\mathcal{E}_{i}|_{*}^{2} + \mathbb{E}|\sum_{i=k}^{\ell}\mathcal{A}^{\ell-i}\mathcal{D}_{i}|_{*}^{2}\right\}, \ \ell = k, \dots, N.$$

If necessary we choose a bound h^0 on the stepsize such that $4 h L_0 < \frac{1}{2}$ holds for all $h < h^0$ and conclude that

$$\mathbb{E}|\mathcal{E}_{\ell}|_{*}^{2} \leq 8 \mathbb{E}|\mathcal{E}_{k-1}|_{*}^{2} + 8 \mathbb{E}|\sum_{i=k}^{\ell} \mathcal{A}^{\ell-i} \mathcal{D}_{i}|_{*}^{2} + 8 L T \frac{1}{N} \sum_{i=k-1}^{\ell-1} \mathbb{E}|\mathcal{E}_{i}|_{*}^{2}.$$

We now apply Lemma 2.11 with $a_{\ell} := 0$, $\ell = 1, \ldots, k - 2$ and $a_{\ell} := \mathbb{E}|\mathcal{E}_{\ell}|^2_*$, $\ell = k - 1, \ldots, N$, and obtain the intermediate result

$$\max_{\ell=k-1,...,N} \mathbb{E}|\mathcal{E}_{\ell}|_{*}^{2} \leq \hat{S} \Big\{ \mathbb{E}|\mathcal{E}_{k-1}|_{*}^{2} + \max_{\ell=k,...,N} \mathbb{E}|\sum_{i=k}^{\ell} \mathcal{A}^{\ell-i} \mathcal{D}_{i}|_{*}^{2} \Big\}, \ \hat{S} := 8 \exp(8LT) \ . \ (A.5)$$

It remains to deal with the term labelled 4), i.e. the perturbations D_i in \mathcal{D}_i . We decompose D_i , and, analogously, \mathcal{D}_i into

 $\begin{aligned} D_i &= R_i + S_i = R_i + \sum_{j=1}^k S_{j,i-j+1}, & \mathcal{D}_i = \mathcal{R}_i + \mathcal{S}_i = \mathcal{R}_i + \sum_{j=1}^k \mathcal{S}_{j,i-j+1}, \\ \text{where } S_{j,i-j+1} \text{ is } \mathcal{F}_{t_{i-j+1}} \text{-measurable with } \mathbb{E}(S_{j,i-j+1} | \mathcal{F}_{t_{i-j}}) = 0 \text{ for } i = k, \dots, N \\ \text{and } j &= 1, \dots, k. \text{ Then } \mathbb{E}\langle \mathcal{A}^{\ell_1} \mathcal{S}_{j_1,i_1}, \mathcal{A}^{\ell_2} \mathcal{S}_{j_2,i_2} \rangle_* = 0 \text{ for } i_1 \neq i_2, \text{ and by similar} \\ \text{computations as above we obtain} \end{aligned}$

$$\mathbb{E}|\sum_{i=k}^{\ell} \mathcal{A}^{\ell-i} \mathcal{D}_{i}|_{*}^{2} \leq 2 \ (\ell-k+1) \sum_{i=k}^{\ell} \mathbb{E}|\mathcal{A}^{\ell-i} \mathcal{R}_{i}|_{*}^{2} + 2 \ k \ \sum_{i=k}^{\ell} \sum_{j=1}^{k} \mathbb{E}|\mathcal{A}^{\ell-i} \mathcal{S}_{j,i-j+1}|_{*}^{2}$$
$$\leq 2 \sum_{i=k}^{\ell} \left(\frac{T}{h} \ \mathbb{E}|\mathcal{R}_{i}|_{*}^{2} + k \sum_{j=1}^{k} \mathbb{E}|\mathcal{S}_{j,i-j+1}|_{*}^{2}\right).$$

Inserting this into the intermediate result (A.5) we obtain

$$\max_{\ell=k-1,...,N} \mathbb{E}|\mathcal{E}_{\ell}|_{*}^{2} \leq \hat{S}\Big\{\mathbb{E}|\mathcal{E}_{k-1}|_{*}^{2} + 2\sum_{i=k}^{\ell} \Big(\frac{T}{h} \mathbb{E}|\mathcal{R}_{i}|_{*}^{2} + k\sum_{j=1}^{k} \mathbb{E}|\mathcal{S}_{j,i-j+1}|_{*}^{2}\Big)\Big\},\$$

and

$$\max_{\ell=k-1,\dots,N} \mathbb{E}|E_{\ell}|^{2} \leq c^{*} \hat{S} \Big\{ c_{*} \max_{\ell=0,\dots,k-1} \mathbb{E}|E_{\ell}|^{2} + 2 c_{11}^{*} \max_{\ell=k,\dots,N} \Big(\frac{T^{2}}{h^{2}} \mathbb{E}|R_{\ell}|^{2} + \frac{kT}{h} \sum_{j=1}^{k} \mathbb{E}|S_{j,\ell-j+1}|^{2} \Big) \Big\}.$$

Taking the square root yields the final estimate

$$\begin{aligned} &\max_{\ell=k-1,\dots,N} \|E_{\ell}\|_{L_{2}} \\ &\leq \sqrt{c^{*}} \hat{S} \bigg\{ \sqrt{c_{*}} \max_{\ell=0,\dots,k-1} \|E_{\ell}\|_{L_{2}} + \sqrt{2c_{11}^{*}} \max_{\ell=1,\dots,N} \bigg(\frac{T}{h} \|R_{\ell}\|_{L_{2}} + \sqrt{\frac{kT}{h}} \sum_{j=1}^{k} \|S_{j,\ell-j+1}\|_{L_{2}}^{2} \bigg) \bigg\} \\ &\leq S \bigg\{ \max_{\ell=0,\dots,k-1} \|E_{\ell}\|_{L_{2}} + \max_{\ell=1,\dots,N} \bigg(\frac{\|R_{\ell}\|_{L_{2}}}{h} + \frac{\sqrt{\sum_{j=1}^{k} \|S_{j,\ell-j+1}\|_{L_{2}}^{2}}}{\sqrt{h}} \bigg) \bigg\}, \end{aligned}$$

which completes the proof. \square

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