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GENERAL ORDER CONDITIONS FOR STOCHASTIC PARTITIONED RUNGE–KUTTA METHODS

SVERRE ANMARKRUD, KRISTIAN DEBRABANT, AND ANNE KVÆRNØ

ABSTRACT. In this paper stochastic partitioned Runge–Kutta (SPRK) methods are considered. A general order theory for SPRK methods based on stochastic B-series and multicolored, multishaped rooted trees is developed. The theory is applied to prove the order of some known methods, and it is shown how the number of order conditions can be reduced in some special cases, especially that the conditions for preserving quadratic invariants can be used as simplifying assumptions.

1. INTRODUCTION

Whenever the right hand side of an ordinary differential equation (ODE) can be split into different parts with different features, it might be worth trying to solve the different parts by different methods. Distinguished examples are stiff/nonstiff parts solved by implicit/explicit methods, fast/slow parts solved by multirate methods, or partitioned symplectic methods for separable Hamiltonian systems. The latter seems to have been the aim when partitioned methods for stochastic differential equations (SDEs) are constructed, e.g. [23, 6, 18, 14, 19, 13, 29]. Although higher order methods have been constructed, to our knowledge there is still no complete order theory for stochastic partitioned Runge–Kutta (SPRK) methods. Such an order theory based on multicolored, multishaped rooted trees is the main contribution in this paper. The theory is valid for both Itô and Stratonovich SDEs.

Constructing high order methods for SDEs is a notoriously nontrivial task. One problem is the required high order stochastic increments. Another is the huge number of order conditions that have to be fulfilled. For this reason, most methods are constructed for problems with some special structure, like linear SDEs, SDEs with additive noise, or separable SDEs. We will demonstrate how the order theory can be simplified in some of these cases. In particular, Sanz-Serna and Abia [27, 1] have proven that for Runge–Kutta and partitioned Runge–Kutta methods preserving quadratic invariants only order conditions related to rootless trees have to be satisfied. A similar result has been proved for stochastic Runge–Kutta methods in [3]. This theory depends on the product rule valid for Stratonovich integrals and is therefore only applicable for Stratonovich SDEs. It will be extended to Stratonovich SPRK methods here.

In this paper we consider a system of stochastic partitioned differential equations with Q partitions and M diffusion terms,

$$(1) \quad X^{(q)}(t) = x_0^{(q)} + \sum_{m=0}^M \int_0^t g_m^{(q)}(X^{(1)}(s), X^{(2)}(s), \dots, X^{(Q)}(s)) \star dW_m(s), \quad q = 1, \dots, Q,$$

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for which we will also use the abbreviated form

$$(2) \quad dX^{(q)}(t) = \sum_{m=0}^M g_m^{(q)}(X^{(1)}(t), X^{(2)}(t), \dots, X^{(Q)}(t)) \star dW_m(t), \quad X^{(q)}(0) = x_0^{(q)}.$$

To simplify the notation the deterministic terms are represented by $m = 0$, such that $dW_0(s) = ds$, while W_m , $m = 1, \dots, M$, denote one-dimensional and pairwise independent Wiener processes. The integrals w.r.t. the Wiener processes are interpreted as either Itô integrals, $\star dW_m(s) = dW_m(s)$, or Stratonovich integrals, $\star dW_m(s) = \circ dW_m(s)$. We also define the vector of initial values, $x_0 = [x_0^{(1)}, x_0^{(2)}, \dots, x_0^{(Q)}]$. Furthermore, we assume that the coefficients $g_m^{(q)} : \mathbb{R}^{d_1} \times \dots \times \mathbb{R}^{d_Q} \rightarrow \mathbb{R}^{d_q}$ are sufficiently smooth, and that the conditions of the existence and uniqueness theorem [24] are satisfied. The systems are considered to be autonomous. Nonautonomous systems can be included by extending one of the variables $X^{(q)}$, $q = 1, \dots, Q$, by t , or by considering the equation $t' = 1$ as a separate partitioning.

Denote the numerical approximation of $X^{(q)}$ at time t_n by $Y_n^{(q)}$. For the solution of (1) we consider general, s -stage SPRK methods given by

$$(3a) \quad H_i^{(q)} = Y_n^{(q)} + \sum_{m=0}^M \sum_{j=1}^s Z_{i,j}^{(q,m)} g_m^{(q)}(H_j^{(1)}, \dots, H_j^{(Q)}), \quad i = 1, \dots, s,$$

$$(3b) \quad Y_{n+1}^{(q)} = Y_n^{(q)} + \sum_{m=0}^M \sum_{i=1}^s \gamma_i^{(q,m)} g_m^{(q)}(H_i^{(1)}, \dots, H_i^{(Q)}),$$

for partitions $q = 1, 2, \dots, Q$. The coefficients $\gamma_i^{(q,m)}$ and $Z_{i,j}^{(q,m)}$, $i, j = 1, \dots, s$, include random variables that depend on the stepsize h . The random variables used in each step are assumed to be i.i.d., and also h might change from step to step. To simplify the notation, here and in the following we omit to indicate this by an additional index n on h , the random variables used, the coefficients $\gamma^{(q,m)} = (\gamma_i^{(q,m)})_{i=1}^s$ and $Z^{(q,m)} = (Z_{i,j}^{(q,m)})_{i,j=1}^s$, and the stage values $H_i^{(q)}$. This will in particular also hold for the Wiener increments $\Delta W_m = \int_{t_n}^{t_n+h} dW_m(s)$.

The coefficients of an SPRK method can be gathered in a generalized Butcher tableau. In the frequently encountered case that there exist matrix functions $Z^{(q)}$ and vector functions $\gamma^{(q)}$ such that with some vectors of random variables ξ_m it holds that $Z^{(q,m)} = Z^{(q)}(\xi_m)$ and $\gamma^{(q,m)} = \gamma^{(q)}(\xi_m)$ for $m = 1, \dots, M$ and $q = 1, \dots, Q$, we will write the Butcher tableau as follows:

$$\begin{array}{c|c} Z^{(1,0)} & Z^{(1,m)} \\ \hline \vdots & \vdots \\ \hline Z^{(Q,0)} & Z^{(Q,m)} \\ \hline \hline (\gamma^{(1,0)})^\top & (\gamma^{(1,m)})^\top \\ \hline \vdots & \vdots \\ \hline (\gamma^{(Q,0)})^\top & (\gamma^{(Q,m)})^\top \end{array}, \quad m = 1, \dots, M.$$

Remark 1. *The splitting (1) is sometimes called a horizontal splitting. The results are equally valid for a vertical splitting, that is SDEs split by*

$$(4) \quad \hat{X}(t) = \hat{x}_0 + \sum_{q=1}^Q \sum_{m=0}^M \int_0^t \hat{g}_m^{(q)}(\hat{X}(s)) \star dW_m(s), \quad q = 1, \dots, Q,$$

as (1) can, assuming $d_1 = \dots = d_Q$ (extending the $X^{(q)}$ by zero components were necessary), be transformed to (4) by

$$\hat{X}(t) = \sum_{q=1}^Q X^{(q)}(t), \quad g_m^{(q)}(X^{(1)}(t), X^{(2)}(t), \dots, X^{(Q)}(t)) = \hat{g}_m^{(q)} \left(\sum_{q=1}^Q X^{(q)}(t) \right).$$

By summing up the expressions of (3) the corresponding SPRK method becomes

$$(5a) \quad \hat{H}_i = \hat{Y}_n + \sum_{q=1}^Q \sum_{m=0}^M \sum_{j=1}^s Z_{i,j}^{(q,m)} \hat{g}_m^{(q)}(\hat{H}_j), \quad i = 1, \dots, s,$$

$$(5b) \quad \hat{Y}_{n+1} = \hat{Y}_n + \sum_{q=1}^Q \sum_{m=0}^M \sum_{i=1}^s \gamma_i^{(q,m)} \hat{g}_m^{(q)}(\hat{H}_i),$$

where \hat{Y}_n is the numerical approximation of $\hat{X}(t_n)$.

Example 1. Consider the Langevin equation of motion

$$(6a) \quad \dot{r} = v,$$

$$(6b) \quad \dot{v} = f(r, t) - \alpha v + \beta(t),$$

which describes the evolution of a particle with unit mass, coordinate $r(t)$ and velocity $v(t)$. The particle is affected by three forces: $f(r, t)$, a friction force αv for a friction coefficient $\alpha \geq 0$ and thermal white noise $\beta(t)$. Representing the noise term as $\beta W(t)$ with some constant β , this can be written as a proper SDE, split vertically

$$(7) \quad \begin{pmatrix} dR(t) \\ dV(t) \end{pmatrix} = \underbrace{\begin{pmatrix} V(t) \\ -\alpha V(t) \end{pmatrix}}_{\hat{g}_0^{(1)}} dt + \underbrace{\begin{pmatrix} 0 \\ f(R(t), t) \end{pmatrix}}_{\hat{g}_0^{(2)}} dt + \underbrace{\begin{pmatrix} 0 \\ \beta \end{pmatrix}}_{\hat{g}_1^{(2)}} dW,$$

or horizontally by

$$\begin{aligned} \begin{pmatrix} dR(t) \\ dU_1(t) \end{pmatrix} &= \underbrace{\begin{pmatrix} U_1(t) + U_2(t) \\ -\alpha(U_1(t) + U_2(t)) \end{pmatrix}}_{g_0^{(1)}} dt, \\ dU_2(t) &= \underbrace{f(R(t), t)}_{g_0^{(2)}} dt + \underbrace{\beta}_{g_1^{(2)}} dW(t) \end{aligned}$$

with $V(t) = U_1(t) + U_2(t)$.

Grønbech-Jensen and Farago [11] proposed the following scheme to solve the Langevin equation,

$$(8a) \quad R_{n+1} = R_n + bhV_n + \frac{bh^2}{2} f(R_n, t_n) + \frac{bh}{2} \beta \Delta W,$$

$$(8b) \quad V_{n+1} = aV_n + \frac{h}{2} (af(R_n, t_n) + f(R_{n+1}, t_{n+1})) + b\beta \Delta W,$$

with

$$a = \frac{1 - \frac{\alpha h}{2}}{1 + \frac{\alpha h}{2}}, \quad b = \frac{1}{1 + \frac{\alpha h}{2}}.$$

This can be reformulated equivalently as an SPRK method with $\hat{H}_1 = \begin{pmatrix} R_n \\ V_n \end{pmatrix}$ and

$$\begin{aligned} \begin{pmatrix} \hat{H}_{2,1} \\ \hat{H}_{2,2} \end{pmatrix} &= \begin{pmatrix} R_n \\ V_n \end{pmatrix} + \frac{h}{2} \begin{pmatrix} \hat{H}_{1,2} + \hat{H}_{2,2} \\ -\alpha(\hat{H}_{1,2} + \hat{H}_{2,2}) \end{pmatrix} + h \begin{pmatrix} 0 \\ f(\hat{H}_{1,1}, t_n) \end{pmatrix} + \Delta W \begin{pmatrix} 0 \\ \beta \end{pmatrix}, \\ \begin{pmatrix} R_{n+1} \\ V_{n+1} \end{pmatrix} &= \begin{pmatrix} R_n \\ V_n \end{pmatrix} + \frac{h}{2} \begin{pmatrix} \hat{H}_{1,2} + \hat{H}_{2,2} \\ -\alpha(\hat{H}_{1,2} + \hat{H}_{2,2}) \end{pmatrix} + \frac{h}{2} \begin{pmatrix} 0 \\ f(\hat{H}_{1,1}, t_n) + f(\hat{H}_{2,1}, t_{n+1}) \end{pmatrix} \\ &\quad + \Delta W \begin{pmatrix} 0 \\ \beta \end{pmatrix}, \end{aligned}$$

where $\hat{H}_{i,j}$ denotes component j of \hat{H}_i . Due to the special structure of the problem, the scheme (8) can be interpreted as an application of several SPRK methods (with $M = 1$), of which we mention the following two:

$$(9) \quad \begin{array}{c} \begin{array}{cc|cc} 0 & 0 & 0 & 0 \\ \hline \frac{h}{2} & \frac{h}{2} & 0 & 0 \\ \hline 0 & 0 & 0 & 0 \\ h & 0 & \Delta W_m & 0 \\ \hline \frac{h}{2} & \frac{h}{2} & 0 & 0 \\ \hline \frac{h}{2} & \frac{h}{2} & \Delta W_m & 0 \end{array} , \quad \begin{array}{c} \begin{array}{cc|cc} 0 & 0 & 0 & 0 \\ \hline \frac{h}{2} & \frac{h}{2} & \frac{\Delta W_m}{2} & \frac{\Delta W_m}{2} \\ \hline 0 & 0 & 0 & 0 \\ h & 0 & \Delta W_m & 0 \\ \hline \frac{h}{2} & \frac{h}{2} & \frac{\Delta W_m}{2} & \frac{\Delta W_m}{2} \\ \hline \frac{h}{2} & \frac{h}{2} & \frac{\Delta W_m}{2} & \frac{\Delta W_m}{2} \end{array} , \quad m = 1, \dots, M. \end{array} \end{array}$$

While the method to the left might be the obvious choice, it is only convergent for partitioned problems for which there is no noise in the first partitioning. In this particular case it coincides with the method to the right. In Section 3.1 we will prove this method to be of strong order 1 and weak order 2 when applied to more general partitioned SDEs with additive noise.

In Section 2, a general order theory for SPRKs is developed. The theory is based on stochastic B-series and multicolored, multishaped rooted trees. In Section 3 two particular cases are studied: SDEs with additive noise and separable problems. We will here present examples of how the B-series theory can be used to find the order of some given methods. In Section 4, it is shown how the number of order conditions can be further reduced in the case of methods preserving quadratic invariants of Stratonovich SDEs.

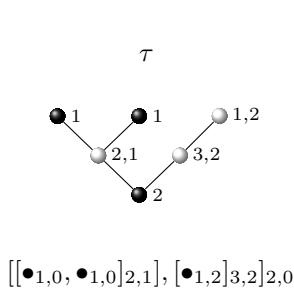
2. ORDER THEORY

B-series for deterministic ODEs were introduced by J. C. Butcher [8] in 1963. B-series for SDEs were developed by Burrage and Burrage [4, 5, 7] for strong convergence of Stratonovich SDEs, by Komori, Mitsui and Sugiura [17] and Komori [16] for weak convergence of Stratonovich SDEs, and by Rößler [25, 26] for weak convergence in both the Itô and Stratonovich case. A unified theory for B-series encompassing both weak and strong convergence for both Itô and Stratonovich SDEs was given in [9]. In the following we will generalize this to SPRKs.

Our first goal is to find B-series representations of (3), and we begin by assuming $X^{(q)}(h)$ can be written as a B-series $B^{(q)}(\phi, x_0; h)$,

$$B^{(q)}(\phi, x_0; h) = \sum_{\tau \in T_q} \alpha(\tau) \cdot \phi(\tau)(h) \cdot F(\tau)(x_0),$$

where T_q is the set of shaped, colored, rooted trees as defined below. The terms $\alpha(\tau)$ are combinatoric terms. The elementary weight functions $\phi(\tau)(h)$ are stochastic integrals or random variables, and $F(\tau)(x_0)$ are the elementary differentials. To simplify the presentation, we assume that all



$$\alpha(\tau) = \frac{1}{2}$$

$$F(\tau) = (D_{2,3}g_0^{(2)})(D_{11}g_1^{(2)})(g_0^{(1)}, g_0^{(1)}), (D_1g_2^{(3)})g_2^{(1)}$$

$$\phi(\tau) = \int_0^h \left(\int_0^s s_1^2 dW_1(s_1) \right) \left(\int_0^s W_2(s_2) dW_2(s_2) \right) ds$$

$$\Phi(\tau) = \sum_{i,j,k,l,m,n=1}^s \gamma_i^{(2,0)} Z_{i,j}^{(2,1)} Z_{j,k}^{(1,0)} Z_{j,l}^{(1,0)} Z_{i,m}^{(3,2)} Z_{m,n}^{(1,2)}$$

$$\rho(\tau) = \frac{9}{2}$$

FIGURE 1. An example of a shaped, colored, rooted tree and its corresponding functions.

elementary differentials exist and all considered B-series converge. Otherwise, one has to consider truncated B-series and discuss the remainder term [25].

Definition 1 (Trees and combinatorial coefficients). *The set of shaped, rooted trees*

$$T = T_1 \cup T_2 \cup \dots \cup T_Q$$

where

$$T_q = \{\emptyset_q\} \cup T_{q,0} \cup T_{q,1} \cup \dots \cup T_{q,M}$$

for $q = 1, \dots, Q$ is recursively defined as follows:

(i) The graph $\bullet_{q,m}$ with only one vertex of shape q and color m belongs to $T_{q,m}$.

(ii) If $\tau_1, \tau_2, \dots, \tau_\kappa \in T \setminus \{\emptyset_1, \dots, \emptyset_Q\}$, then $[\tau_1, \tau_2, \dots, \tau_\kappa]_{q,m} \in T_{q,m}$, where $[\tau_1, \tau_2, \dots, \tau_\kappa]_{q,m}$ denotes the tree formed by joining the subtrees $\tau_1, \tau_2, \dots, \tau_\kappa$ each by a single branch to a common root of shape q and color m .

Further, we define $\alpha(\tau)$ as

$$\alpha(\emptyset_q) = 1, \alpha(\bullet_{q,m}) = 1, \alpha(\tau = [\tau_1, \dots, \tau_\kappa]_{q,m}) = \frac{1}{r_1! r_2! \dots r_\kappa!} \prod_{k=1}^{\kappa} \alpha(\tau_k),$$

where $r_1, r_2, \dots, r_\kappa$ count equal trees among $\tau_1, \tau_2, \dots, \tau_\kappa$.

Remark 2. Rooted trees can be represented in the bracket notation, as used in the definition, or illustrated as graphs (see Fig. 1). To ease the reading, deterministic nodes are in the latter represented as black nodes, with the $m = 0$ omitted, while stochastic nodes are white. Figure 1 also gives the corresponding examples for the functions of trees that will be defined in the following.

Definition 2 (Elementary differentials). For a tree $\tau \in T$ the elementary differential is a mapping $F(\tau): \mathbb{R}^{d_1} \times \dots \times \mathbb{R}^{d_Q} \rightarrow \mathbb{R}^d$ defined recursively by

(i) $F(\emptyset_q)(x_0) = x_0^{(q)}, \emptyset_q \in T_q,$

(ii) $F(\bullet_{q,m})(x_0) = g_m^{(q)}(x_0),$

(iii) If $\tau = [\tau_1, \tau_2, \dots, \tau_\kappa]_{q,m} \in T_{q,m}$, then

$$F(\tau)(x_0) = (D_{q_1 \dots q_\kappa} g_m^{(q)})(x_0)(F(\tau_1)(x_0), F(\tau_2)(x_0), \dots, F(\tau_\kappa)(x_0))$$

where q_k is the shape of τ_k , $k = 1, \dots, \kappa$, and $D_{q_1 \dots q_\kappa} = \frac{\partial^\kappa}{\partial x^{q_1} \dots \partial x^{q_\kappa}}$ denotes the derivative operator of order κ .

Fundamental for this work is the following lemma which says that if $Y^{(q)}(h)$ can be written as a B-series, then $f(Y^{(1)}(h), \dots, Y^{(Q)}(h))$ can also be written as a B-series. This is a trivial extension of the lemma found in [9].

Lemma 1. *If $Y^{(q)}(h) = B^{(q)}(\phi, x_0; h)$, $q = 1, \dots, Q$, are some B-series and $f \in C^\infty(\mathbb{R}^{d_1} \times \dots \times \mathbb{R}^{d_Q}, \mathbb{R}^d)$, then $f(Y^{(1)}(h), \dots, Y^{(Q)}(h))$ can be written as a formal series of the form*

$$(10) \quad f(Y^{(1)}(h), \dots, Y^{(Q)}(h)) = \sum_{u \in U_f} \beta(u) \cdot \psi_\phi(u)(h) \cdot G(u)(x_0),$$

where

(i) U_f is a set of trees derived from T , by $\bullet_f \in U_f$, and if $\tau_1, \tau_2, \dots, \tau_\kappa \in T \setminus \{\emptyset_1, \dots, \emptyset_Q\}$, then $[\tau_1, \tau_2, \dots, \tau_\kappa]_f \in U_f$.

(ii) $G(\bullet_f)(x_0) = f(x_0)$ and

$$G([\tau_1, \tau_2, \dots, \tau_\kappa]_f)(x_0) = (D_{q_1, \dots, q_\kappa} f)(x_0)(F(\tau_1)(x_0), \dots, F(\tau_\kappa)(x_0)).$$

(iii) $\beta(\bullet_f) = 1$ and

$$\beta([\tau_1, \tau_2, \dots, \tau_\kappa]_f) = \frac{1}{r_1! r_2! \dots r_q!} \prod_{k=1}^{\kappa} \alpha(\tau_k),$$

where r_1, r_2, \dots, r_q count equal trees among $\tau_1, \tau_2, \dots, \tau_\kappa$.

(iv) $\psi_\phi(\bullet_f)(h) = 1$ and $\psi_\phi([\tau_1, \tau_2, \dots, \tau_\kappa]_f)(h) = \prod_{k=1}^{\kappa} \phi(\tau_k)(h)$.

If we apply Lemma 1 to the functions $g_m^{(q)}$ in (1) we get

$$g_m^{(q)}(X^{(1)}(h), \dots, X^{(Q)}(h)) = \sum_{u \in U_{g_m^{(q)}}} \beta(u) \cdot \psi_\phi(u)(h) \cdot G(u)(x_0).$$

By the definitions of trees, $T_{q,m}$, and elementary differentials, $F(\tau)(x_0)$, we can write this as

$$(11) \quad g_m^{(q)}(X^{(1)}(h), \dots, X^{(Q)}(h)) = \sum_{\tau \in T_{q,m}} \alpha(\tau) \cdot \phi'_{q,m}(\tau)(h) \cdot F(\tau)(x_0),$$

where

$$\phi'_{q,m}(\tau)(h) = \begin{cases} 1, & \text{if } \tau = \bullet_{q,m}, \\ \prod_{k=1}^{\kappa} \phi(\tau_k)(h), & \text{if } \tau = [\tau_1, \dots, \tau_\kappa]_{q,m} \in T_{q,m}. \end{cases}$$

We now write the exact solutions of (1) as B-series and use (11) to obtain

$$\sum_{\tau \in T_q} \alpha(\tau) \cdot \phi(\tau)(h) \cdot F(\tau)(x_0) = x_0^{(q)} + \sum_{m=0}^M \sum_{\tau \in T_{q,m}} \alpha(\tau) \cdot \int_0^h \phi'_{q,m}(\tau)(s) \star dW_m(s) F(\tau)(x_0).$$

Comparing term by term we see that

$$\phi(\emptyset_q)(h) = 1, \text{ and } \phi(\tau)(h) = \int_0^h \phi'_{q,m}(\tau)(s) \star dW_m(s) \text{ for all } \tau \in T_{q,m}, m = 0, 1, \dots, M, q = 1, \dots, Q.$$

With induction on the height of τ we have proven the following theorem.

Theorem 1. *The exact solutions $X^{(q)}(h)$ of (1), $q = 1, \dots, Q$, can be written as B-series $B^{(q)}(\phi, x_0; h)$ with*

$$\begin{aligned} \phi(\emptyset_q)(h) &= 1, & \phi(\bullet_{q,m})(h) &= \Delta W_m(h), \\ \phi([\tau_1, \tau_2, \dots, \tau_\kappa]_{q,m})(h) &= \int_0^h \prod_{j=1}^{\kappa} \phi(\tau_j(s)) \star dW_m(s), \\ & \text{for all } [\tau_1, \tau_2, \dots, \tau_\kappa]_{q,m} \in T_{q,m}, \quad q = 1, \dots, Q, \quad m = 0, 1, \dots, M. \end{aligned}$$

A similar result can be found for the numerical solution of (1) by the s -stage SPRK method (3).

Theorem 2. *The numerical solutions $Y_1^{(q)}$ as well as the stage values can be written in terms of B-series*

$$H_i^{(q)} = B^{(q)}(\Psi_i, x_0; h), \quad Y_1^{(q)} = B^{(q)}(\Phi, x_0; h)$$

for all $i = 1, \dots, s$, $q = 1, \dots, Q$, with

$$(12a) \quad \Psi_i(\emptyset_q)(h) = 1, \quad \Psi_i(\bullet_{q,m})(h) = \sum_{j=1}^s Z_{i,j}^{(q,m)},$$

$$(12b) \quad \Psi_i([\tau_1, \dots, \tau_\kappa]_{q,m})(h) = \sum_{j=1}^s Z_{i,j}^{(q,m)} \prod_{k=1}^{\kappa} \Psi_j(\tau_k)(h)$$

and

$$(13a) \quad \Phi(\emptyset_q)(h) = 1, \quad \Phi(\bullet_{q,m})(h) = \sum_{i=1}^s \gamma_i^{(q,m)},$$

$$(13b) \quad \Phi(\tau = [\tau_1, \dots, \tau_\kappa]_{q,m})(h) = \sum_{i=1}^s \gamma_i^{(q,m)} \prod_{k=1}^{\kappa} \Psi_i(\tau_k)(h).$$

Proof. Write $H_i^{(q)}$ as B-series,

$$H_i^{(q)} = \sum_{\tau \in T_q} \alpha(\tau) \Psi_i(\tau)(h) F(\tau)(x_0),$$

for $i = 1, \dots, s$, $q = 1, \dots, Q$, where as usual the product of vectors is understood componentwise. Inserted into (3a) and using (11) this gives

$$H_i^{(q)} = x_0^{(q)} + \sum_{m=0}^M \sum_{j=1}^s \sum_{\tau \in T_{q,m}} \alpha(\tau) \left(Z_{i,j}^{(q,m)} \cdot \Psi'_{j,q,m}(\tau)(h) \right) F(\tau)(x_0).$$

A term by term comparison yields (12). The proof of (13) is similar. \square

The local order of accuracy of the SPRK method can now be decided by comparing the B-series of the exact and the numerical solution. Applying a time transformation, it is hereby sufficient to consider B-series expansions around $t = 0$, as done in Theorems 1 and 2. First, we need to define the tree order.

Definition 3. *The order of a tree $\tau \in T$ is defined by*

$$\rho(\emptyset_q) = 0, \quad \rho(\tau = [\tau_1, \dots, \tau_\kappa]_{q,m}) = \sum_{k=1}^{\kappa} \rho(\tau_k) + \begin{cases} 1 & \text{for } m = 0, \\ \frac{1}{2} & \text{otherwise.} \end{cases}$$

The following theorem relates the global order of accuracy to the local order. Here, we assume that method (3) is constructed such that $\Phi(\tau)(h) = \mathcal{O}(h^{\rho(\tau)})$ for all $\tau \in T$ for mean square convergence, respectively $E \prod_{k=1}^{\kappa} \Phi(\tau_k) = \mathcal{O}(h^{\sum_{k=1}^{\kappa} \rho(\tau_k)})$ for all $\tau_1, \dots, \tau_{\kappa} \in T$, $\kappa \in \mathbb{N}$, for weak convergence.

Theorem 3. *The method has mean square global order p if*

$$(14a) \quad \Phi(\tau)(h) = \phi(\tau)(h) + \mathcal{O}(h^{p+\frac{1}{2}}), \quad \forall \tau \in T \text{ with } \rho(\tau) \leq p,$$

$$(14b) \quad E\Phi(\tau)(h) = E\phi(\tau)(h) + \mathcal{O}(h^{p+1}), \quad \forall \tau \in T \text{ with } \rho(\tau) \leq p + \frac{1}{2},$$

and weak consistency of order p if and only if

$$(15) \quad E \prod_{k=1}^{\kappa} \Phi(\tau_k) = E \prod_{k=1}^{\kappa} \phi(\tau_k) + \mathcal{O}(h^{p+1}) \quad \text{whenever} \quad \sum_{k=1}^{\kappa} \rho(\tau_k) \leq p + \frac{1}{2}.$$

Here, the $\mathcal{O}(\cdot)$ -notation refers to $h \rightarrow 0$ and, especially in (14a), to the L^2 -norm. The result (14) was first proved in [5], while (15) is a consequence of a result of Milstein [22], see [9] for details. A list of all trees with $\rho(\tau) \leq 2$ and their corresponding functions are given in Appendix A.

3. TWO SPECIAL CASES

The amount of order conditions to be satisfied for higher order methods is quite overwhelming. E. g. for a method of strong order 1.5 with $Q = 2$ and scalar noise, order conditions for 122 different trees need to be satisfied (some of them being trivially fulfilled, though).

The virtue of partitioned methods becomes clear when applied to problems for which some underlying structure can be exploited, in the sense that many elementary differentials will be zero, and the corresponding trees can thus be ignored in the B-series. The main task is to identify those trees. In this subsection, the idea will be demonstrated with two distinguished examples from Milstein et. al [20, 23, 21], SDEs with additive noise and separable systems, both with two partitionings.

3.1. SDEs with additive noise. We consider partitioned problems with $Q = 2$ with additive noise, thus

$$(16a) \quad dX^{(1)} = g_0^{(1)}(X^{(1)}, X^{(2)}, t)dt + \sum_{m=1}^M g_m^{(1)}(t)dW_m(t),$$

$$(16b) \quad dX^{(2)} = g_0^{(2)}(X^{(1)}, X^{(2)}, t)dt + \sum_{m=1}^M g_m^{(2)}(t)dW_m(t).$$

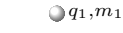
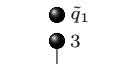
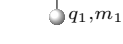
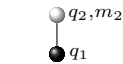
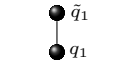
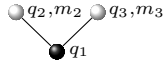
To deal with the nonautonomous case, let us include a third partition:

$$dX^{(3)} = 1dt.$$

The problem structure induces that many elementary differentials will be zero, and the corresponding trees can thus be ignored in the B-series. Restricting to nonvanishing elementary differentials, we only need to consider trees characterized by the following properties:

- There are no nodes $\bullet_{3,m}$ for $m \neq 0$, since the corresponding $g_m^{(3)} = 0$.
- Nodes $\bullet_{3,0}$ have no branches, since $g_0^{(3)}$ is constant.
- The stochastic nodes $\bullet_{q,m}$ with $q = 1, 2$ and $m \neq 0$ can only be followed by one or more nodes $\bullet_{3,0}$, since $g_m^{(q)}$ is only time-dependent.

TABLE 1. Relevant trees up to order two and corresponding functions for two-partitioned SDEs with additive noise (16). Here, $q_1, q_2, q_3 \in \{1, 2\}$, while $\tilde{q}_1 \in \{1, 2, 3\}$. The weights $\Phi(\tau)$ correspond to method (9) extended by $Z^{(3,0)} = Z^{(2,0)}$ and $\gamma^{(3,0)} = \gamma^{(2,0)}$.

No	τ	$\rho(\tau)$	$\phi(\tau)$	$\Phi(\tau)$
1		$\frac{1}{2}$	ΔW_{m_1}	ΔW_{m_1}
2		1	h	h
6		$\frac{3}{2}$	$\int_0^h s dW_{m_1}(s)$	$\frac{1}{2}h\Delta W_{m_1}$
7		$\frac{3}{2}$	$\int_0^h W_{m_2}(s) ds$	$\frac{1}{2}h\Delta W_{m_2}$
8		2	$\frac{1}{2}h^2$	$\frac{1}{2}h^2$
9		2	$\int_0^h W_{m_2}(s)W_{m_3}(s) ds$	$\frac{1}{2}h\Delta W_{m_2}\Delta W_{m_3}$

Of the trees listed in the appendix, this leaves us with the trees no. 1, 2, 6, 7, 8 and 9 and even those can be simplified since $q_1 \neq 3$ for the trees no. 6, 7, 8 and 9, and $q_2 = 3$ for tree no. 6. The trees together with their corresponding functions are listed in Table 1.

Let the SDE (16) be solved by method (9). For the third partition, we use $Z^{(3,0)} = Z^{(2,0)}$ and $\gamma^{(3,0)} = \gamma^{(2,0)}$. The corresponding weights $\Phi(\tau)$ are given in Table 1. From Theorem 3 we can conclude that the method is of strong order 1, as the order conditions for the trees no. 6 and 7 are only fulfilled in expectation, and that it is of weak order 2.

The Langevin equation (6) of Example 1 is an example of an SDE (16), thus we can conclude that the method proposed in [11] is of strong order 1 and weak order 2. A further example of an SDE (16) is the following:

Example 2 (Stochastic version of the Jansen and Rit Neural Mass Model). *Let X_0, X_1 and X_2 describe the mean postsynaptic potentials of the main neural population, the excitatory and inhibitory interneurons, respectively. The following model has been proposed in [2]:*

$$\begin{aligned}
dX_i(t) &= X_{i+3}(t), \quad i \in \{0, 1, 2\}, \\
dX_3(t) &= [Aa [\mu_3(t) + \text{Sigm}(X_1(t) - X_2(t))] - 2aX_3(t) - a^2X_0(t)] + \sigma_3(t)dW_3(t), \\
dX_4(t) &= [Aa [\mu_4(t) + C_2 \text{Sigm}(C_1X_0(t))] - 2aX_4(t) - a^2X_1(t)] + \sigma_4(t)dW_4(t), \\
dX_5(t) &= [Bb [\mu_5(t) + C_4 \text{Sigm}(C_3X_0(t))] - 2bX_5(t) - b^2X_2(t)] + \sigma_5(t)dW_5(t),
\end{aligned}$$

where $\text{Sigm}(x) = \frac{\nu_{\max}}{1+e^r(\nu_0-x)}$, μ_i and σ_i , $i \in \{3, 4, 5\}$, describe external input and the scaling of the stochastic components, respectively, and A, B, a, b, C_i , $i \in \{1, 2, 3, 4\}$, r and ν_0, ν_{\max} are parameters.

3.2. Separable systems. Consider the following system, which typically can arise from separable Hamiltonian systems, discussed e.g. in [21, 4.2.2]:

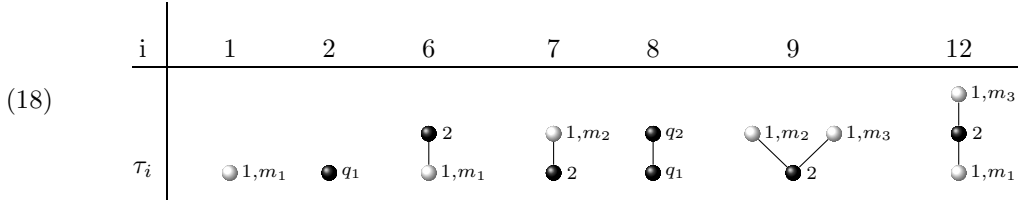
$$(17a) \quad dX^{(1)}(t) = g_0^{(1)}(X^{(2)}(t))dt + \sum_{m=1}^M g_m^{(1)}(X^{(2)}(t)) \star dW_m(t),$$

$$(17b) \quad dX^{(2)}(t) = g_0^{(2)}(X^{(1)}(t))dt.$$

The elementary differentials corresponding to the following trees vanish:

- All trees for which a node of shape q ($q \in \{1, 2\}$) is followed by a node of the same shape, as $\frac{\partial g_m^{(q)}}{\partial x^{(q)}} = 0$ for $m = 0, \dots, M$.
- All trees with nodes $\bullet_{2,m}$ with $m \neq 0$: There is no noise in the second partition, so $g_m^{(2)} = 0$ for $m = 1, \dots, M$.

The remaining trees τ with $\rho(\tau) \leq 2$ are then (for $q_1, q_2 \in \{1, 2\}$ with $q_1 \neq q_2$)



Consider the following method, proposed in [23]:

$\frac{h}{4}$	0	$\frac{3J_{(m,0)}}{2h} - \frac{\Delta W_m}{2}$	0
$\frac{h}{4}$	$\frac{3h}{4}$	$\frac{3J_{(m,0)}}{2h} - \frac{\Delta W_m}{2}$	$-\frac{3J_{(m,0)}}{2h} + \frac{3\Delta W_m}{2}$
0	0		
$\frac{2h}{3}$	0		
$\frac{h}{4}$	$\frac{3h}{4}$	$\frac{3J_{(m,0)}}{2h} - \frac{\Delta W_m}{2}$	$-\frac{3J_{(m,0)}}{2h} + \frac{3\Delta W_m}{2}$
$\frac{2h}{3}$	$\frac{h}{3}$		

where $J_{(m,0)} = \int_{t_n}^{t_n+h} (W_m(s) - W_m(t_n))ds$. It is straightforward to show that $\phi(\tau) = \Phi(\tau)$ for trees $\tau_1, \tau_2, \tau_7, \tau_8$ in (18). Using the fact (see e.g. [15, 10]) that $J_{(m,0)} + J_{(0,m)} = h\Delta W_m$, where $J_{(0,m)} = \int_{t_n}^{t_n+h} (s - t_n) \star dW_m(s)$, we obtain for τ_6 that (remember that for the consistency analysis it is enough to consider the first step of the method, so $t_n = t_0 = 0$)

$$\begin{aligned} \Phi(\tau_6) &= \sum_{i=1}^2 \gamma_i^{(1,m_1)} \sum_{j=1}^2 Z_{i,j}^{(2,0)} = \gamma_2^{(1,m_1)} \sum_{j=1}^2 Z_{2,j}^{(2,0)} \\ &= \left(-\frac{3J_{(m_1,0)}}{2h} + \frac{3\Delta W_{m_1}}{2} \right) \frac{2h}{3} = -J_{(m_1,0)} + \Delta W_{m_1} h \\ &= J_{(0,m_1)} = \phi(\tau_6). \end{aligned}$$

For the remaining two trees we get:

$$\begin{aligned}\Phi(\tau_9) &= \frac{1}{2} \left(\frac{3}{h} J_{(m_2,0)} J_{(m_3,0)} - \Delta W_{m_2} J_{(m_3,0)} - \Delta W_{m_3} J_{(m_2,0)} + h \Delta W_{m_2} \Delta W_{m_3} \right), \\ \Phi(\tau_{12}) &= \frac{1}{2} \left(-\frac{3}{h} J_{(m_1,0)} J_{(m_3,0)} + 3 \Delta W_{m_1} J_{(m_3,0)} + \Delta W_{m_3} J_{(m_1,0)} - h \Delta W_{m_1} \Delta W_{m_3} \right).\end{aligned}$$

From Theorem 3 we can conclude that the method is (remarkably both for Stratonovich and Itô SDEs) of strong order 1.5, since

$$E\Phi(\tau_9) = E\phi(\tau_9) = \begin{cases} \frac{h^2}{2} & \text{if } m_2 = m_3 \\ 0 & \text{otherwise} \end{cases}, \quad E\Phi(\tau_{12}) = E\phi(\tau_{12}) = 0,$$

which follows from the following relations:

$$E\Delta W_m^2 = h, \quad E\Delta W_m J_{(m,0)} = \frac{h^2}{2}, \quad EJ_{(m,0)}^2 = \frac{h^3}{3}.$$

This is of course in accordance with the order result given (for the Stratonovich case) in [20, Theorem 4.3].

4. QUADRATIC INVARIANTS AS SIMPLIFYING ASSUMPTIONS, ROOTLESS TREES

In the previous section, we demonstrated how the special structure of the SDEs can be exploited, in the sense that the elementary differentials corresponding to certain trees are zero, and the corresponding conditions for these trees do not need to be fulfilled. In this section, we will see how for separable Stratonovich SDEs the algebraic relation between the method coefficients that implies the method to preserve quadratic invariants creates certain equivalence classes of trees in the sense that only one condition for each class has to be fulfilled. This is a generalization of a result obtained for deterministic partitioned ODEs by Abia and Sanz-Serna [1]. It is similar to the nonpartitioned case, which for ODEs is discussed by Sanz-Serna and Abia [27] and generalized to SDEs in Anmarkrud and Kværnø [3]. The conditions for preserving quadratic invariants for non-separable equations with two partitionings have been developed by Hong, Xu and Wang [19], see also Ma and Ding [19], both in the context of symplectic methods. Using similar ideas, it is possible to find conditions for a more general partitioning. However, the extra freedom gained by more partitionings is to some extent lost by the huge number of extra conditions that have to be satisfied. In this paper, the discussion is restricted to separable equations with the number of partitions of (1) to be two, i.e.

$$(20a) \quad dX^{(1)}(t) = \sum_{m=0}^M g_m^{(1)}(X^{(2)}(t)) \circ dW_m(t),$$

$$(20b) \quad dX^{(2)}(t) = \sum_{m=0}^M g_m^{(2)}(X^{(1)}(t)) \circ dW_m(t),$$

which is also the stochastic counterpart to the system discussed in [1].

We assume that the system has a quadratic invariant $I(X^{(1)}, X^{(2)}) = X^{(1)\top} D X^{(2)}$ for a matrix D of the appropriate dimension and arbitrary initial values $X^{(1)}(0), X^{(2)}(0)$.

Example 3 (Synchrotron oscillations). *Consider a stochastically perturbed Hamiltonian system of a pendulum ([28], where we set $\lambda_A = 0$) with Hamiltonians*

$$(21) \quad H_0(p, x, t) = \frac{p^2}{2}, \quad H_1(p, x, t) = \omega^2 \sin(x) \lambda_{Ph}.$$

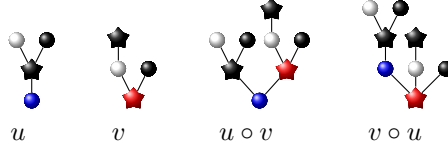


FIGURE 2. Demonstration of the Butcher product (with real colors and shapes instead of indices).

The resulting SDE system

$$(22) \quad dp = -\omega^2 \sin(x)dt - \lambda_{Ph}\omega^2 \cos(x) \circ dW(t), \quad dx = pdt$$

preserves symplecticity, i. e. a quadratic invariant [21, Section 4.1].

As in [12, chapter IV.2.2] one can prove the following theorem:

Theorem 4. *If the coefficients of the SPRK method (3) satisfy*

$$(23) \quad \gamma_i^{(1,m_1)} \gamma_j^{(2,m_2)} = \gamma_j^{(2,m_2)} Z_{j,i}^{(1,m_1)} + \gamma_i^{(1,m_1)} Z_{i,j}^{(2,m_2)} \\ \forall i, j = 1, \dots, s \text{ and } \forall m_1, m_2 = 0, 1, \dots, M$$

then it preserves quadratic invariants of (20) of the form $I(X^{(1)}, X^{(2)}) = X^{(1)\top} DX^{(2)}$.

When the conditions (23) are fulfilled, the number of order conditions stated in Theorem 3 can be reduced. The key ingredient in the analysis is the Butcher product of two trees $u, v \in \tau_1, \tau_2, \dots, \tau_\kappa \in T \setminus \{\emptyset_1, \dots, \emptyset_Q\}$. If

$$(24) \quad u = [u_1, \dots, u_{\kappa_1}]_{q_1, m_1}, \quad v = [v_1, \dots, v_{\kappa_2}]_{q_2, m_2},$$

the Butcher product is defined by

$$u \circ v = [u_1, \dots, u_{\kappa_1}, v]_{q_1, m_1},$$

see Fig. 2. We can now state the following lemma:

Lemma 2. *For Stratonovich SDEs of the form (20) we have for all $u \in T_{q_1, m_1}$ and $v \in T_{q_2, m_2}$*

$$\phi(u)(h) \cdot \phi(v)(h) = \phi(u \circ v)(h) + \phi(v \circ u)(h)$$

for all $q_1, q_2 = 1, 2$ and $m_1, m_2 = 0, 1, \dots, M$.

Proof. As for u and v given by (24) it holds

$$\phi(u)(t) = \int_0^t \prod_{k_1=1}^{\kappa_1} \phi(u_{k_1})(s) \circ dW_{m_1}(s), \quad \phi(v)(t) = \int_0^t \prod_{k_2=1}^{\kappa_2} \phi(v_{k_2})(s) \circ dW_{m_2}(s),$$

the product rule for Stratonovich integrals gives

$$\begin{aligned} \phi(u)(t) \cdot \phi(v)(t) &= \int_0^t \phi(v)(s) \prod_{k_1=1}^{\kappa_1} \phi(u_{k_1})(s) \circ dW_{m_1}(s) \\ &\quad + \int_0^t \phi(u)(s) \prod_{k_2=1}^{\kappa_2} \phi(v_{k_2})(s) \circ dW_{m_2}(s) \\ &= \phi(u \circ v)(t) + \phi(v \circ u)(t). \end{aligned}$$

□

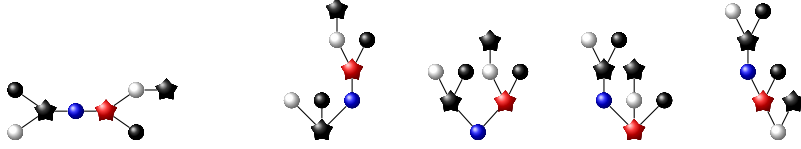


FIGURE 3. An unrooted tree $\hat{\tau}$ to the left, and some of the trees in $\widehat{TS}(\hat{\tau})$. The figure also illustrates the process described in the proof of Theorem 5.

For $q_1 \neq q_2$ this rule also holds for the weight functions of the numerical solution as stated in the following lemma:

Lemma 3. *If an SPRK method of the form (3) satisfies the conditions in Theorem 4, then for all $u \in T_{1,m_1}$ and $v \in T_{2,m_2}$*

$$\Phi(u)(h) \cdot \Phi(v)(h) = \Phi(u \circ v)(h) + \Phi(v \circ u)(h)$$

for all $m_1, m_2 = 0, 1, \dots, m$.

Proof. Multiply both sides of (23) in Theorem 4 by

$$\prod_{k_1=1}^{\kappa_1} \Psi_i(u_{k_1})(h) \prod_{k_2=1}^{\kappa_2} \Psi_j(v_{k_2})(h)$$

and sum over all $i, j = 1, \dots, s$,

$$\begin{aligned} \left(\sum_i \gamma_i^{(1,m_1)} \prod_{k_1=1}^{\kappa_1} \Psi_i(u_{k_1})(h) \right) \left(\sum_j \gamma_j^{(2,m_2)} \prod_{k_2=1}^{\kappa_2} \Psi_j(v_{k_2})(h) \right) = \\ \sum_{i,j} \left(\gamma_j^{(2,m_2)} Z_{j,i}^{(1,m_1)} + \gamma_i^{(1,m_1)} Z_{i,j}^{(2,m_2)} \right) \prod_{k_1=1}^{\kappa_1} \Psi_i(u_{k_1})(h) \prod_{k_2=1}^{\kappa_2} \Psi_j(v_{k_2})(h). \end{aligned}$$

Using (12) and (13) completes the proof. \square

Let TS be the set of trees for which a node of one shape will never be directly followed by a node of the same shape. For separable SDEs, all elementary differentials of trees in $T \setminus TS$ vanish. Given a $\tau \in TS$, let $\hat{\tau}$ be the corresponding unrooted tree, and let $\widehat{TS}(\hat{\tau}) \subset TS$ be the set of trees obtained from $\hat{\tau}$ by assigning one of the nodes as the root, see Fig. 3 for an illustration.

Theorem 5. *Assume that (23) is satisfied. Let $\hat{\tau} \in \widehat{TS}$ be an unrooted tree of order $q \leq p$. If $\phi(\tau)(h) = \Phi(\tau)(h) + \mathcal{O}(h^{p+\frac{1}{2}})$ for one rooted tree $\tau \in \widehat{TS}(\hat{\tau})$ and all rooted trees of order less than q , then it holds that $\phi(\tau)(h) = \Phi(\tau)(h) + \mathcal{O}(h^{p+\frac{1}{2}})$ for all $\tau \in \widehat{TS}(\hat{\tau})$. Similarly, if $E\phi(\tau)(h) = E\Phi(\tau)(h) + \mathcal{O}(h^{p+\frac{1}{2}})$ for one rooted tree $\tau \in \widehat{TS}(\hat{\tau})$ and $\phi(\tau)(h) = \Phi(\tau)(h) + \mathcal{O}(h^{p+\frac{1}{2}})$ for all rooted trees of order less than q , then it holds that $E\phi(\tau)(h) = E\Phi(\tau)(h) + \mathcal{O}(h^{p+\frac{1}{2}})$ for all $\tau \in \widehat{TS}(\hat{\tau})$.*

Proof. The argument is the same as for the non-partitioned case proven in [3]. For trees with one node the theorem is trivially true. Let $\hat{\tau}$ be an unrooted tree of order q and two or more nodes, and let τ be a corresponding rooted tree $\tau \in \widehat{TS}(\hat{\tau})$. Pick one branch v from the root of τ and let the remaining part of τ be u , so that $\tau = u \circ v$. Then u and v have roots of different shapes, and Lemma 3 applies. Clearly, the orders of u and v are less than the order of τ , and by Lemma 2 and the assumptions of the theorem we then have

$$\phi(v \circ u)(h) = \Phi(v \circ u)(h) + \mathcal{O}(h^{p+\frac{1}{2}})$$

respectively




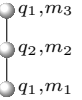

$$E\phi(v \circ u)(h) = E\Phi(v \circ u)(h) + \mathcal{O}(h^{p+\frac{1}{2}}).$$

Because the choice of branch v was arbitrary, this means that this condition is satisfied for all trees with the same graph as τ , but with a root shifted to one of its neighboring nodes. A repeated use of this argument proves the result. The process is illustrated in Fig. 3. \square

Example 4. A simple example for an SPRK method fulfilling (23) is the stochastic Störmer-Verlet method [14], given by the following tableau:

0	0	0	0
$\frac{h}{2}$	$\frac{h}{2}$	$\frac{\Delta W_m}{2}$	$\frac{\Delta W_m}{2}$
$\frac{h}{2}$	0	$\frac{\Delta W_m}{2}$	0
$\frac{h}{2}$	0	$\frac{\Delta W_m}{2}$	0
$\frac{h}{2}$	$\frac{h}{2}$	$\frac{\Delta W_m}{2}$	$\frac{\Delta W_m}{2}$
$\frac{h}{2}$	$\frac{h}{2}$	$\frac{\Delta W_m}{2}$	$\frac{\Delta W_m}{2}$

To analyze its order, note that due to (20) being separable, the elementary differentials corresponding to trees for which a node of shape q ($q \in \{1, 2\}$) is followed by a node of the same shape vanish. Taking into account Theorem 5, the remaining trees τ with $\rho(\tau) \leq \frac{3}{2}$ to be considered are then (for $q_1, q_2 \in \{1, 2\}$ with $q_1 \neq q_2$)

i	1	2	3	5	7
τ_i					

It is straightforward to show that $\phi(\tau) = \Phi(\tau)$ for trees τ_1, τ_2 in (25), and that this is true for τ_3 only if $m_1 = m_2$, while it holds also for $m_1 \neq m_2$ that $E\phi(\tau_3) = E\Phi(\tau_3)$. Thus, for SDEs (20) with multidimensional noise, the above method is only of order $\frac{1}{2}$. For scalar noise, it holds also $\phi(\tau_7) = \Phi(\tau_7)$ and $E\phi(\tau_5) = E\Phi(\tau_5)$, but not $\phi(\tau_5) = \Phi(\tau_5)$, so in this case the method is of order 1. Similarly, it follows that the method is of weak order one for multidimensional noise.

The content of this section demonstrates the potential of the B-series formulation, in the sense that theory developed and well understood for ODEs here directly can be applied for SDEs. It also demonstrates how to use the particular structure of the problem at hand to reduce the set of order conditions that have to be satisfied. This is demonstrated in Table 2. But even if the number of conditions is reduced in this case, the conditions in Theorem 4 are still so restrictive that it is far from a trivial task to construct higher order methods.

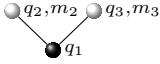
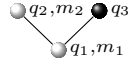

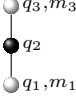

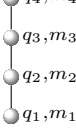
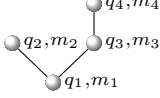
5. CONCLUSION

In this paper we have developed a general B-series theory for SPRK methods. Such methods are rarely applied to general SDEs, they are more likely to be constructed for SDEs with certain structure. We have therefore emphasized how the general theory, summarized in Theorem 3, can be modified to cover a few common examples of such cases. We hope other researchers can find the theory useful for constructing new methods, or to prove order results for applying existing methods to broader classes than originally constructed for.

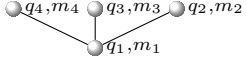
TABLE 2. Total number of trees up to order 2 for SDEs with scalar noise and two partitionings. Here, column (all) gives the number when assuming no restrictions, column (s.p.) the number counting only trees with nonzero elementary differentials for separable problems, and column (q.i.) the number of independent conditions if additionally the quadratic invariant condition is satisfied.

ρ	all	s.p.	q.i.
0.5	2	2	2
1	6	4	3
1.5	22	8	4
2	92	20	9
Sum	122	34	18

No.	τ	$\rho(\tau)$	$\frac{\phi(\tau)}{\Phi(\tau)}$	$F(\tau)$
1		0.5	$\frac{W_{m_1}(h)}{\sum_i \gamma_i^{(q_1, m_1)}}$	$g_{m_1}^{(q_1)}$
2		1	$\frac{h}{\sum_i \gamma_i^{(q_1, 0)}}$	$g_0^{(q_1)}$
3		1	$\frac{\int_0^h W_{m_2}(s) \star dW_{m_1}(s)}{\sum_{i,j} \gamma_i^{(q_1, m_1)} Z_{i,j}^{(q_2, m_2)}}$	$(D_{q_2} g_{m_1}^{(q_1)}) g_{m_2}^{(q_2)}$
4		1.5	$\frac{\int_0^h W_{m_2}(s) W_{m_3}(s) \star dW_{m_1}(s)}{\sum_{i,j,k} \gamma_i^{(q_1, m_1)} Z_{i,j}^{(q_2, m_2)} Z_{i,k}^{(q_3, m_3)}}$	$(D_{q_2 q_3} g_{m_1}^{(q_1)}) (g_{m_2}^{(q_2)}, g_{m_3}^{(q_3)})$
5		1.5	$\frac{\int_0^h \int_0^s W_{m_3}(s_1) \star dW_{m_2}(s_1) \star dW_{m_1}(s)}{\sum_{i,j,k} \gamma_i^{(q_1, m_1)} Z_{i,j}^{(q_2, m_2)} Z_{j,k}^{(q_3, m_3)}}$	$(D_{q_2} g_{m_1}^{(q_1)}) (D_{q_3} g_{m_2}^{(q_2)}) g_{m_3}^{(q_3)}$
6		1.5	$\frac{\int_0^h s \star dW_{m_1}(s)}{\sum_{i,j} \gamma_i^{(q_1, m_1)} Z_{i,j}^{(q_2, 0)}}$	$(D_{q_2} g_{m_1}^{(q_1)}) g_0^{(q_2)}$
7		1.5	$\frac{\int_0^h W_{m_2}(s) ds}{\sum_{i,j} \gamma_i^{(q_1, 0)} Z_{i,j}^{(q_2, m_2)}}$	$(D_{q_2} g_0^{(q_1)}) g_{m_2}^{(q_2)}$
8		2	$\frac{\frac{1}{2} h^2}{\sum_{i,j} \gamma_i^{(q_1, 0)} Z_{i,j}^{(q_2, 0)}}$	$(D_{q_2} g_0^{(q_1)}) g_0^{(q_2)}$

9		2	$\frac{\int_0^h W_{m_2}(s)W_{m_3}(s)ds}{\sum_{i,j,k} \gamma_i^{(q_1,0)} Z_{i,j}^{(q_2,m_2)} Z_{i,k}^{(q_3,m_3)}}$	$(D_{q_2 q_3} g_0^{(q_1)})(g_{m_2}^{(q_2)}, g_{m_3}^{(q_3)})$
10		2	$\frac{\int_0^h W_{m_2}(s)s \star dW_{m_1}(s)}{\sum_{i,j,k} \gamma_i^{(q_1,m_1)} Z_{i,j}^{(q_2,m_2)} Z_{i,k}^{(q_3,0)}}$	$(D_{q_2 q_3} g_{m_1}^{(q_1)})(g_{m_2}^{(q_2)}, g_0^{(q_3)})$
11		2	$\frac{\int_0^h \int_0^s W_{m_3}(s_1) \star dW_{m_2}(s_1) \star dW_{m_1}(s)}{\sum_{i,j,k} \gamma_i^{(q_1,0)} Z_{i,j}^{(q_2,m_2)} Z_{j,k}^{(q_3,m_3)}}$	$(D_{q_2} g_0^{(q_1)})(D_{q_2} g_{m_2}^{(q_2)}, g_{m_3}^{(q_3)})$
12		2	$\frac{\int_0^h \int_0^s W_{m_3}(s_1) ds_1 \star dW_{m_1}(s)}{\sum_{i,j,k} \gamma_i^{(q_1,m_1)} Z_{i,j}^{(q_2,0)} Z_{j,k}^{(q_3,m_3)}}$	$(D_{q_2} g_{m_1}^{(q_1)})(D_{q_3} g_0^{(q_2)}, g_{m_3}^{(q_3)})$
13		2	$\frac{\int_0^h \int_0^s s_1 \star dW_{m_2}(s_1) \star dW_{m_1}(s)}{\sum_{i,j,k} \gamma_i^{(q_1,m_1)} Z_{i,j}^{(q_2,m_2)} Z_{j,k}^{(q_3,0)}}$	$(D_{q_2} g_{m_1}^{(q_1)})(D_{q_3} g_{m_2}^{(q_2)}, g_0^{(q_3)})$
14		2	$\frac{\int_0^h \int_0^s \int_0^{s_1} dW_{m_4}(s_2) \star dW_{m_3}(s_2) \star dW_{m_2}(s_1) \star dW_{m_1}(s)}{\sum_{i,j,k,l} \gamma_i^{(q_1,m_1)} Z_{i,j}^{(q_2,m_2)} Z_{j,k}^{(q_3,m_3)} Z_{k,l}^{(q_4,m_4)}}$	$(D_{q_2} g_{m_1}^{(q_1)})(D_{q_3} g_{m_2}^{(q_2)})(D_{q_4} g_{m_3}^{(q_3)}) g_{m_4}^{(q_4)}$
15		2	$\frac{\int_0^h W_{m_2}(s) \int_0^s W_{m_4}(s_1) \circ dW_{m_3}(s_1) \circ dW_{m_1}(s)}{\sum_{i,j,k,l} \gamma_i^{(q_1,m_1)} Z_{i,j}^{(q_2,m_2)} Z_{i,k}^{(q_3,m_3)} Z_{k,l}^{(q_4,m_4)}}$	$(D_{q_2 q_3} g_{m_1}^{(q_1)})(g_{m_2}^{(q_2)}, (D_{q_4} g_{m_3}^{(q_3)}) g_{m_4}^{(q_4)})$

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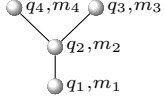


2

$$-\frac{\int_0^h W_{m_4}(s)W_{m_3}(s)W_{m_2}(s) \circ dW_{m_1}(s)}{\sum_{i,j,k,l} \gamma_i^{(q_1, m_1)} Z_{i,j}^{(q_2, m_2)} Z_{i,k}^{(q_3, m_3)} Z_{i,l}^{(q_4, m_4)}}$$

$$(D_{q_2 q_3 q_4} g_{m_1}^{(q_1)})(g_{m_2}^{(q_2)}, g_{m_3}^{(q_3)}, g_{m_4}^{(q_4)})$$

17



2

$$-\frac{\int_0^h \int_0^s W_{m_4}(s_1)W_{m_3}(s_1) \circ dW_{m_2}(s_1) \circ dW_{m_1}(s)}{\sum_{i,j,k,l} \gamma_i^{(q_1, m_1)} Z_{i,j}^{(q_2, m_2)} Z_{j,k}^{(q_3, m_3)} Z_{j,l}^{(q_4, m_4)}}$$

$$(D_{q_2} g_{m_1}^{(q_1)})(D_{q_3 q_4} g_{m_2}^{(q_2)})(g_{m_3}^{(q_3)}, g_{m_4}^{(q_4)})$$

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