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# On the stochastic Magnus expansion and its application to SPDEs 

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#### Abstract

We derive a stochastic version of the Magnus expansion for the solution of linear systems of Itô stochastic differential equations (SDEs). The goal of this paper is twofold. First, we prove existence and a representation formula for the logarithm associated to the solution of the matrix-valued SDEs. Second, we propose a new method for the numerical solution of stochastic partial differential equations (SPDEs) based on spatial discretization and application of the stochastic Magnus expansion. A notable feature of the method is that it is fully parallelizable. We also present numerical tests in order to asses the accuracy of the numerical schemes.


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## 1 Introduction

Let $d, q \in \mathbb{N}$ and consider the linear matrix-valued Itô SDE

$$
\left\{\begin{array}{l}
d X_{t}=B_{t} X_{t} d t+A_{t}^{(j)} X_{t} d W_{t}^{j},  \tag{1.1}\\
X_{0}=I_{d},
\end{array}\right.
$$

with $A^{(1)}, \ldots, A^{(q)}, B$ being real $(d \times d)$-matrix-valued bounded stochastic processes, $I_{d}$ the identity $(d \times d)$ matrix and $W=\left(W^{1}, \cdots, W^{q}\right)$ a $q$-dimensional standard Brownian motion. In 1.1), as well as anywhere throughout the paper, we use Einstein summation convention to imply summation of terms containing $W^{j}$, over the index $j$ from 1 to $q$.

In the deterministic case, i.e. $A^{(j)} \equiv 0, j=1, \cdots, q$, 1.1) reduces to the matrix-valued ODE

$$
\left\{\begin{array}{l}
\frac{d}{d t} X_{t}=B_{t} X_{t},  \tag{1.2}\\
X_{0}=I_{d},
\end{array}\right.
$$

[^0]which admits an explicit solution, in terms of matrix exponential, in the time-homogeneous case. Namely, if $B_{t} \equiv B$, the unique solution to 1.2 reads as
$$
X_{t}=e^{t B}, \quad t \geq 0
$$

However, in the non-autonomous case, the ODE 1.2 does not admit an explicit solution. In particular, if $B_{t}$ is not constant, the solution $X_{t}$ typically differs from $e^{\int_{0}^{t} B_{s} d s}$. This is due to the fact that, in general, $B_{t}$ and $B_{s}$ do not commute for $t \neq s$. As it turns out, a representation of the solution in terms of a matrix exponential is still possible, at least for short times, i.e.

$$
\begin{equation*}
X_{t}=e^{Y_{t}} \tag{1.3}
\end{equation*}
$$

for $t \geq 0$ suitably small and $Y_{t}$ real valued $(d \times d)$-matrix. Moreover, $Y$ admits a semi-explicit expansion as a series of iterated integrals involving nested Lie commutators of the function $B$ at different times. Such representation is known as Magnus expansion [20] and its first terms read as

$$
\begin{equation*}
Y_{t}=\int_{0}^{t} B_{s} d s+\frac{1}{2} \int_{0}^{t} d s \int_{0}^{s}\left[B_{s}, B_{u}\right] d u+\frac{1}{6} \int_{0}^{t} d s \int_{0}^{s} d u \int_{0}^{u}\left(\left[B_{s},\left[B_{u}, B_{r}\right]\right]+\left[B_{r},\left[B_{u}, B_{s}\right]\right]\right) d r+\cdots \tag{1.4}
\end{equation*}
$$

The Magnus expansion has a wide range of physical applications and the related literature has grown increasingly over the last decades (see, for instance, the excellent survey paper [3] and the references given therein).

In the stochastic case, when $j=1, B_{t} \equiv 0$ and $A$ is constant, i.e. $A_{t}(\omega) \equiv A$, the Itô equation 1.1) reduces to

$$
\left\{\begin{array}{l}
d X_{t}=A X_{t} d W_{t} \\
X_{0}=I_{d}
\end{array}\right.
$$

whose explicit solution can be easily proved to be of the form (1.3), with

$$
Y_{t}=-\frac{1}{2} A^{2} t+A W_{t}, \quad t \geq 0
$$

In general, when the matrices $A_{t}^{(j)}, A_{s}^{(j)}, B_{t}, B_{s}$ with $t \neq s$ do not commute, an explicit solution to (1.1) is not known. For instance, in the non-commutative case, neither the equation

$$
\left\{\begin{array}{l}
d X_{t}=B d t+A X_{t} d W_{t}  \tag{1.5}\\
X_{0}=I_{d}
\end{array}\right.
$$

nor the equation

$$
\left\{\begin{array}{l}
d X_{t}=A_{t} X_{t} d W_{t} \\
X_{0}=I_{d}
\end{array}\right.
$$

admit an explicit solution, save some particular cases (see for instance the second example in Section 2.3). Among the approximation tools that were developed in the literature to solve stochastic differential equations, including (1.1), some Magnus-type expansions that extend (1.3)-1.4 were derived in different contexts . We now go on to describe our contribution to this stream of literature, and then to firm our results within the existing ones. In particular, a detailed comparison with existing stochastic Magnus expansions previously derived by other authors will be provided below, in the last subsection.

Description of the main results. In this paper we derive a Magnus-type representation formula for the solution to the Itô $\operatorname{SDE}(1.1)$, which is $(1.3)$ together with

$$
\begin{equation*}
Y_{t}=Y_{t}^{(1)}+Y_{t}^{(2)}+Y_{t}^{(3)}+\cdots \quad t \in[0, \tau] \tag{1.6}
\end{equation*}
$$

for $\tau$ suitably small, strictly positive stopping time. In analogy with the deterministic Magnus expansion, the general term $Y^{(n)}$ can be expressed recursively, and contains iterated stochastic integrals of nested Lie commutators of the processes $B, A^{(j)}$ at different times.

In the case $j=1$, the first two terms of the expansion read as

$$
\begin{aligned}
& Y_{t}^{(1)}=\int_{0}^{t} B_{s} d s+\int_{0}^{t} A_{s} d W_{s} \\
& Y_{t}^{(2)}=\frac{1}{2} \int_{0}^{t}\left(\left[B_{s}, \int_{0}^{s} B_{u} d u+\int_{0}^{s} A_{u} d W_{u}\right]-A_{s}^{2}\right) d s+\frac{1}{2} \int_{0}^{t}\left[A_{s}, \int_{0}^{s} B_{u} d u+\int_{0}^{s} A_{u} d W_{u}\right] d W_{s} .
\end{aligned}
$$

For example, in the case of SDE 1.5 the latter can be reduced to

$$
\begin{aligned}
& Y_{t}^{(1)}=B t+A W_{t} \\
& Y_{t}^{(2)}=[A, B]\left(\frac{1}{2} t W_{t}-\int_{0}^{t} W_{s} d s\right)-\frac{1}{2} A^{2} t
\end{aligned}
$$

Notice that the last expressions do not contain stochastic integrals. In fact, in the general autonomous case, and if $j=1$, all the iterated stochastic integrals in $Y^{(n)}$ can be solved for any $n$ (see Corollary 5.2.4 in [13]). Therefore, in this case the expansion becomes numerically computable by only approximating Lebesgue integrals, as opposed to stochastic Runge-Kutta schemes, which typically require the numerical approximation of stochastic integrals. As we shall see in the numerical tests in Section 2.3, this feature allows us to choose a sparser time-grid in order to save computation time. This feature is also preserved in some non-autonomous cases as illustrated in Section 2.3,

A notable feature of the expansion is the possibility of parallelizing the computation of its terms. In contrast to standard iterative methods, which require the solution at a given time-step in order to go through the next step in the iteration, the discretization of the integrals in the terms $Y^{(n)}$ can be done simultaneously for all the time steps. Conclusively, this entails the possibility of parallelizing over all times in the time-grid and makes the numerical implementation of the stochastic Magnus expansion perfectly GPU applicable.

As it often happens when deriving convergent (either asymptotically or absolutely) expansions, a formal derivation precedes the rigorous result: that is what we do for 1.3 - 1.6 in Section 2.2 . Just like the derivation of the deterministic Magnus expansion relies on the possibility of writing the logarithm $Y$ as the solution to an ODE, in the stochastic case the first step consists in representing $Y$ as the solution to an SDE. Such representation of $Y$ will be more involved compared to the deterministic case because of the presence of the second order derivatives of the exponential map coming from the application of Itô's formula. This is a distinctive feature of our derivation with respect to other analogous results in the Stratonovich setting where the standard chain-rule applies. With the SDE representation for $Y$ at hand, the expansion 1.6 stems, like in the deterministic case, from applying a Dyson-type perturbation procedure to the SDE solved by $Y$.

In the deterministic case, the convergence of the Magnus expansion 1.4 to the exact logarithm of the solution to $\sqrt{1.2}$ was studied by several authors, who proved progressively sharper lower bounds on the
maximum $\bar{t}$ such that the convergence to the exact solution is assured for any $t \in[0, \bar{t}]$. At the best of our knowledge, the sharpest estimate was given in 23, namely

$$
\begin{equation*}
\bar{t} \geq \sup \left\{t \geq 0 \mid \int_{0}^{t}\left\|B_{s}\right\| d s<\pi\right\} \tag{1.7}
\end{equation*}
$$

where $\left\|B_{s}\right\|$ denotes the spectral norm. Note that the existence of a real logarithm of $X_{t}$ is an issue that underlies the study of the convergence of the Magnus expansion. We state here our main result, proved in Section 3 which deals with these matters in the stochastic case, when the coefficients $B, A^{(j)}$ in (1.1) are progressively measurable processes. We defer a comparison with previous convergence results for stochastic Magnus-type expansions to the next subsection. We denote by $\mathcal{M}^{d \times d}$ the space of the $(d \times d)$ matrices with real entries. Also, for an $\mathcal{M}^{d \times d}$-valued stochastic process $M=\left(M_{t}\right)_{t \in[0, T]}$, we set $\|M\|_{T}:=$ $\||M|\|_{L^{\infty}([0, T] \times \Omega)}$, where $|\cdot|$ denotes the Frobenius (Euclidean entry-wise) norm.

Theorem 1.1. Let $A^{(1)}, \ldots, A^{(q)}$ and $B$ be bounded, progressively measurable, $\mathcal{M}^{d \times d}$-valued processes defined on a filtered probability space $\left(\Omega, \mathcal{F}, P,\left(\mathcal{F}_{t}\right)_{t \geq 0}\right)$ equipped with a standard $q$-dimensional Brownian motion $W=\left(W^{1}, \cdots, W^{q}\right)$. For $T>0$ let also $X=\left(X_{t}\right)_{t \in[0, T]}$ be the solution to 1.1). There exists a strictly positive stopping time $\tau \leq T$ such that:
(i) $X_{t}$ has a real logarithm $Y_{t} \in \mathcal{M}^{d \times d}$ up to time $\tau$, i.e.

$$
\begin{equation*}
X_{t}=e^{Y_{t}}, \quad 0 \leq t<\tau \tag{1.8}
\end{equation*}
$$

(ii) the following representation holds $P$-almost surely:

$$
\begin{equation*}
Y_{t}=\sum_{n=0}^{\infty} Y_{t}^{(n)}, \quad 0 \leq t<\tau \tag{1.9}
\end{equation*}
$$

where $Y^{(n)}$ is the $n$-th term in the stochastic Magnus expansion as defined in (2.13);
(iii) there exists a positive constant $C$, only dependent on $\left\|A^{(1)}\right\|_{T}, \ldots,\left\|A^{(q)}\right\|_{T},\|B\|_{T}, T$ and $d$, such that

$$
\begin{equation*}
P(\tau \leq t) \leq C t, \quad t \in[0, T] \tag{1.10}
\end{equation*}
$$

The proof of (i) relies on the continuity of $X$ together with a standard representation for the matrix logarithm. The key point in the proof of (ii) consists in showing that $X_{t}^{\varepsilon, \delta}$ and its logarithm $Y_{t}^{\varepsilon, \delta}$ are holomorphic as functions of $(\varepsilon, \delta)$, where $X_{t}^{\varepsilon, \delta}$ represents the solution of 1.1 when $A^{(j)}$ and $B$ are replaced by $\varepsilon A^{(j)}$ and $\delta B$, respectively. Once this is established, the representation $\sqrt{1.9}$ follows from observing that, by construction, the series in 1.9 is exactly the formal power series of $Y_{t}^{\varepsilon, \delta}$ at $(\varepsilon, \delta)=(1,1)$. To prove the holomorphicity of $X_{t}^{\varepsilon, \delta}$ we follow the same approach typically adopted to prove regularity properties of stochastic flows. Namely, in Lemma 3.3 we state some maximal $L^{p}$ and Hölder estimates (with respect to the parameters) for solutions to SDEs with random coefficients and combine them with the Kolmogorov continuity theorem. Finally, the proof of (iii) owes one more time to the $L^{p}$ estimates in Lemma 3.3 and to a Sobolev embedding theorem to obtain pointwise estimates w.r.t. the parameters $(\varepsilon, \delta)$ above.

The last part of the paper, Section 4, is devoted to the application of the stochastic Magnus expansion to the numerical solution of parabolic stochastic partial differential equations (SPDEs). The idea is to discretize the SPDE only in space and then approximate the resulting linear matrix-valued SDE by truncating the series in 1.8 - 1.9 . The goal here is to propose the application of stochastic Magnus expansions as novel approximation tools for SPDEs; we study the error of this approximating procedure only numerically, in a case where an explicit benchmark is available, and we defer the theoretical error analysis to further studies.

Review of the literature and comparison. Stochastic generalizations of the Magnus expansions were proposed by several authors. To the best of our knowledge, we recognize mainly two streams of research.

The beginning of the first one can be traced back to the work [2], where the author derived exponential stochastic Taylor expansions (see also [1], 13] for general stochastic Taylor series) of the solution of a system of Stratonovich SDEs with values on a manifold $\mathcal{M}$, i.e.

$$
\left\{\begin{array}{l}
d X_{t}=B\left(X_{t}\right) d t+A^{(j)}\left(X_{t}\right) \circ d W_{t}^{j}  \tag{1.11}\\
X_{0}=x_{0}
\end{array}\right.
$$

with $B, A^{(j)}$ being smooth, deterministic and autonomous vector fields on $\mathcal{M}$. The stochastic flow of 1.11) is represented in terms of the exponential map of a stochastic vector field $Y$, i.e.

$$
X_{t}\left(x_{0}\right)=\exp Y_{t}\left(x_{0}\right)
$$

the vector field $Y$ being expressed by an infinite series of iterated stochastic integrals multiplying nested commutators of the vector fields $B, A^{(j)}$. This representation is proved up to a strictly positive stopping time and extends some previous results in [11, [27] for the commutative case and in [29], [16] for the nilpotent case. Refinements of [2] were proved in [6] making the expansion of $Y$ more explicit. Later, numerical methods based on these representations were proposed in [7] and [8. Such techniques, known as Castell-Gaines methods, require the approximation of the solution to a time-dependent ODE besides the approximation of iterated stochastic integrals. Truncating the expansion of $Y$ at a specified order, these schemes turn out to be asymptotically efficient in the sense of Newton [24].

If $\mathcal{M}=\mathcal{M}^{d \times d}$ and the vector fields are linear, then (1.11) reduces to the Stratonovich version of 1.1) with $B, A^{(j)}$ constant matrices, and the representation of $X$ given in [2] can be seen as a stochastic Magnus expansion, in that the exponential map of $Y$ reduces to the multiplication by a matrix exponential. In fact, in this case the expansion in [2] becomes explicit in terms of iterated stochastic integrals, and can be shown to coincide with the expansion in this paper by applying Itô-Stratonovich conversion formula. In the very interesting paper [19], the authors study several computational aspects of numerical schemes stemming from the truncated Magnus expansion, in which the iterated stochastic integrals are approximated by their conditional expectation. Besides showing that asymptotic efficiency holds for an arbitrary number of Brownian components, they compare the theoretical accuracy with the one of analogous schemes based on Dyson (or Neumann) series, which are obtained by applying stochastic Taylor expansion directly on the equation. They find that, although the theoretical accuracy of Magnus schemes is not superior, Magnus-based approximations seem more accurate than their Dyson counterparts in practice. They also discuss the computational cost deriving from approximating the iterated stochastic integrals and the matrix exponentiation, in relation to different features of the problem such as the dimension and the number of Brownian motions, as well as to the order of the numerical scheme.

The second stream of literature is explicitly aimed at extending the original Magnus results to stochastic settings and can be traced back to [5] where the Magnus expansion is derived via formal arguments for a linear system of Stratonovich SDEs with deterministic coefficients. Clearly, in the autonomous case such expansion coincides with the one obtained by Ben Arous, whereas in the non-autonomous case, $B \equiv 0$ and $j=1$, it is formally equivalent to the deterministic Magnus expansion with all the Lebesgue integrals replaced by Stratonovich ones. The authors of [5] do not address the convergence of the Magnus expansion,
but rather study computational aspects of the resulting approximation, in particular in comparison with Runge-Kutta stochastic schemes. The authors of 21] consider the Ito SDE (1.1] with constant coefficients, and propose to resolve via Euler method the SDE (2.11) for the logarithm of the solution. In [28] the Magnus expansion for the Stratonovich version of (1.1) with deterministic coefficients is applied to solve non-linear SDEs; however, the error analysis of the truncated expansion seems flawed, since the fact that the Magnus series converges only up to a positive stopping time is completely overlooked.

We now go on to discuss the contribution of this paper with respect to the existing literature. In the first place, Theorem 1.1 on the convergence of the Magnus expansion requires very weak conditions on the coefficients, which are stochastic processes satisfying the sole assumption of progressive measurability. This is a novel aspect compared to the results in [2], [6], which surely cover a wider class of SDEs, but under the assumption of time-independent deterministic coefficients. We point out that this feature is also relevant in light of the fact that our result is stated for Itô SDEs as opposed to Stratonovich ones. Indeed, while this difference might appear as minor in the Markovian case, where a simple conversion formula exists (cf. [10] and [18]), it becomes substantial in the case of progressively measurable coefficients. We also point out that, even in the Markovian non-autonomous case, convergence issues were not discussed in [5] and [19].

Another novel aspect of our result concerns the estimate 1.10 for the cumulative distribution function of the stopping time $\tau$ up to which the Magnus series converges to the real logarithm of the solution: this kind of estimate was unknown even in the autonomous case. Theorem 11 in [2] (see also [6]) provides an asymptotic estimate for the truncation error of the logarithm, which in the linear case studied in this paper would read as

$$
Y_{t}=\sum_{n=1}^{N} Y_{t}^{(n)}+t^{\frac{N+1}{2}} R_{t}, \quad t<T
$$

with $R$ bounded in probability. Although this type of result holds for the general SDE (1.11), it is weaker than Theorem 1.1 in the linear case. In fact, it can be obtained by 1.10 together with the standard estimate $\left\|\sup _{0 \leq s \leq t}\left|Y_{s}^{(n)}\right|\right\|_{L^{2}(\Omega)} \leq C t^{\frac{N+1}{2}}$, but not the other way around.

Finally, we also propose a new method for the numerical solution of SPDEs based on spatial discretization and application of the stochastic Magnus expansion.

The rest of the paper is structured as follows. In Section 2 we derive and test the Magnus expansion. In particular, Section 2.1 contains the definitions of the operators through which the terms $Y^{(n)}$ in 1.8 - 1.9 will be defined and some preliminary results that will be used to derive the expansion. Section 2.2 contains a formal derivation of (1.8)- (1.9), which will be tested numerically on some specific examples in Section 2.3 Section 3 is entirely devoted to the proof of Theorem 1.1. In Section 4 we present the application of the Magnus expansion to the numerical solution of parabolic SPDEs. In particular, in Section 4.1 we recall some general facts about stochastic Cauchy problems, in Section 4.2 we introduce the finite-difference-Magnus approximation scheme, while in Section 4.3 we check the effectiveness of the proposed approach through numerical tests. Finally, Appendix A contains a representations for the first and second order differentials of the matrix exponential map.

## 2 Stochastic Magnus expansion

In this section we define the terms in the expansion 1.9 and present some numerical tests to confirm the accuracy of the approximate solutions to stemming from the truncation of the series.

### 2.1 Preliminaries

Let $\mathcal{M}^{d \times d}$ be the vector space of $(d \times d)$ real-valued matrices. Throughout the paper we denote by $[\cdot, \cdot]$ the standard Lie brackets operation, i.e.

$$
[M, N]:=M N-N M, \quad M, N \in \mathcal{M}^{d \times d}
$$

and by $\|\cdot\|$ the spectral norm on $\mathcal{M}^{d \times d}$. Also, we denote by $\beta_{k}, k \in \mathbb{N}_{0}$, the Bernoulli numbers defined as the derivatives of the function $x \mapsto x /\left(e^{x}-1\right)$ computed at $x=0$. For sake of convenience we report the first three Bernoulli numbers: $\beta_{0}=1, \beta_{1}=-\frac{1}{2}, \beta_{2}=\frac{1}{6}$. Note also that $\beta_{2 m+1}=0$ for any $m \in \mathbb{N}$.

We now define the operators that we will use in the sequel. For a fixed $\Sigma \in \mathcal{M}^{d \times d}$, we let:
$-\operatorname{ad}_{\Sigma}^{j}: \mathcal{M}^{d \times d} \rightarrow \mathcal{M}^{d \times d}$, for $j \in \mathbb{N}_{0}$, be the linear operators defined as

$$
\begin{aligned}
\operatorname{ad}_{\Sigma}^{0}(M) & :=M \\
\operatorname{ad}_{\Sigma}^{j}(M) & :=\left[\Sigma, \operatorname{ad}_{\Sigma}^{j-1}(M)\right], \quad j \in \mathbb{N}
\end{aligned}
$$

To ease notation we also set $\operatorname{ad}_{\Sigma}:=\operatorname{ad}_{\Sigma}^{1}$;

- $e^{\operatorname{ad} \Sigma}: \mathcal{M}^{d \times d} \rightarrow \mathcal{M}^{d \times d}$ be the linear operator defined as

$$
\begin{equation*}
e^{\operatorname{ad}_{\Sigma}}(M):=\sum_{n=0}^{\infty} \frac{1}{n!} \operatorname{ad}_{\Sigma}^{n}(M)=e^{\Sigma} M e^{-\Sigma} \tag{2.1}
\end{equation*}
$$

where $e^{\Sigma}:=\sum_{j=0}^{\infty} \frac{\Sigma^{j}}{j!}$ is the standard matrix exponential;

- $\mathcal{L}_{\Sigma}: \mathcal{M}^{d \times d} \rightarrow \mathcal{M}^{d \times d}$ be the linear operator defined as

$$
\begin{equation*}
\mathcal{L}_{\Sigma}(M):=\int_{0}^{1} e^{\operatorname{ad}_{\tau \Sigma}}(M) d \tau=\sum_{n=0}^{\infty} \frac{1}{(n+1)!} \operatorname{ad}_{\Sigma}^{n}(M) \tag{2.2}
\end{equation*}
$$

$-Q_{\Sigma}: \mathcal{M}^{d \times d} \times \mathcal{M}^{d \times d} \rightarrow \mathcal{M}^{d \times d}$ be the bi-linear operator defined as

$$
\begin{align*}
\mathcal{Q}_{\Sigma}(M, N) & :=\mathcal{L}_{\Sigma}(M) \mathcal{L}_{\Sigma}(N)+\int_{0}^{1} \tau\left[\mathcal{L}_{\tau \Sigma}(N), e^{\operatorname{ad}_{\tau \Sigma}}(M)\right] d \tau  \tag{2.3}\\
& =\sum_{n=0}^{\infty} \sum_{m=0}^{\infty} \frac{\operatorname{ad}_{\Sigma}^{n}(M)}{(n+1)!} \frac{\operatorname{ad}_{\Sigma}^{m}(N)}{(m+1)!}+\sum_{n=0}^{\infty} \sum_{m=0}^{\infty} \frac{\left[\operatorname{ad}_{\Sigma}^{n}(N), \operatorname{ad}_{\Sigma}^{m}(M)\right]}{(n+m+2)(n+1)!m!} \tag{2.4}
\end{align*}
$$

Proposition 2.1 (Itô formula). Let $Y$ be an $\mathcal{M}^{d \times d}$-valued Itô process of the form

$$
\begin{equation*}
d Y_{t}=\mu_{t} d t+\sigma_{t}^{j} d W_{t}^{j} \tag{2.5}
\end{equation*}
$$

Then we have

$$
d e^{Y_{t}}=\left(\mathcal{L}_{Y_{t}}\left(\mu_{t}\right)+\frac{1}{2} \sum_{j=1}^{q} \mathcal{Q}_{Y_{t}}\left(\sigma_{t}^{j}, \sigma_{t}^{j}\right)\right) e^{Y_{t}} d t+\mathcal{L}_{Y_{t}}\left(\sigma_{t}^{j}\right) e^{Y_{t}} d W_{t}^{j}
$$

Proof. The statement follows from the multi-dimensional Itô formula (see, for instance, [25]) combined with Lemma A. 1 and applied to the exponential process $e^{Y_{t}}$.

We also have the following inversion formula for the operator $\mathcal{L}_{\Sigma}$.
Lemma 2.2 (Baker, 1905). Let $\Sigma \in \mathcal{M}^{d \times d}$. The operator $\mathcal{L}_{\Sigma}$ is invertible if and only if the eigenvalues of the linear operator $a d_{\Sigma}$ are different from $2 m \pi, m \in \mathbb{Z} \backslash\{0\}$. Furthermore, if $\|\Sigma\|<\pi$, then

$$
\begin{equation*}
\mathcal{L}_{\Sigma}^{-1}(M)=\sum_{k=0}^{\infty} \frac{\beta_{k}}{k!} a d_{\Sigma}^{k}(M), \quad M \in \mathcal{M}^{d \times d} \tag{2.6}
\end{equation*}
$$

For a proof to Lemma 2.2 we refer the reader to [3].

### 2.2 Formal derivation

In this section we perform formal computations to derive the terms $Y^{(n)}$ appearing in the Magnus expansion (1.9). Although such computations are heuristic at this stage, they are meant to provide the reader with an intuitive understanding of the principles that underlie the expansion procedure. Their validity will be proved a fortiori, in Section 3, in order to prove Theorem 1.1.

Let $\left(\Omega, \mathcal{F}, P,\left(\mathcal{F}_{t}\right)_{t \geq 0}\right)$ a filtered probability space. Assume that, for any $\varepsilon, \delta \in \mathbb{R}$, the process $X^{\varepsilon, \delta}=$ $\left(X_{t}^{\varepsilon, \delta}\right)_{t \geq 0}$ solves the Itô SDE

$$
\left\{\begin{array}{l}
d X_{t}^{\varepsilon, \delta}=\delta B_{t} X_{t}^{\varepsilon, \delta} d t+\varepsilon A_{t}^{(j)} X_{t}^{\varepsilon, \delta} d W_{t}^{j}  \tag{2.7}\\
X_{0}^{\varepsilon, \delta}=I_{d}
\end{array}\right.
$$

and that it admits the exponential representation

$$
\begin{equation*}
X_{t}^{\varepsilon, \delta}=e^{Y_{t}^{\varepsilon, \delta}} \tag{2.8}
\end{equation*}
$$

with $Y^{\varepsilon, \delta}$ being an $\mathcal{M}^{d \times d}$-valued Itô process. Clearly, if $(\varepsilon, \delta)=(1,1)$, then 2.7 - (2.8) reduce to (1.1)-(1.3). of the form 2.5 . Then, Proposition 2.1 yields

$$
\begin{align*}
\varepsilon A_{t}^{(j)} & =\mathcal{L}_{Y_{t}^{\varepsilon, \delta}}\left(\sigma_{t}^{j}\right), \quad j=1, \ldots, q  \tag{2.9}\\
\delta B_{t} & =\mathcal{L}_{Y_{t}^{\varepsilon, \delta}}\left(\mu_{t}\right)+\frac{1}{2} \sum_{j=1}^{q} \mathcal{Q}_{Y_{t}^{\varepsilon, \delta}}\left(\sigma_{t}^{j}, \sigma_{t}^{j}\right) \tag{2.10}
\end{align*}
$$

Inverting now 2.9-2.10, in accord with 2.6), one obtains

$$
\begin{aligned}
\sigma_{t}^{j} & =\mathcal{L}_{Y_{t}^{\varepsilon, \delta}}^{-1}\left(\varepsilon A_{t}^{(j)}\right)=\varepsilon \sum_{k=0}^{\infty} \frac{\beta_{k}}{k!} \operatorname{ad}_{Y_{t}^{\varepsilon, \delta}}^{k}\left(A_{t}^{(j)}\right), \quad j=1, \ldots, q \\
\mu_{t} & =\mathcal{L}_{Y_{t}^{\varepsilon, \delta}}^{-1}\left(\delta B_{t}-\frac{1}{2} \sum_{j=1}^{q} \mathcal{Q}_{Y_{t}^{\varepsilon, \delta}}\left(\sigma_{t}^{j}, \sigma_{t}^{j}\right)\right) \\
& =\sum_{k=0}^{\infty} \frac{\beta_{k}}{k!} \operatorname{ad}_{Y_{t}^{\varepsilon, \delta}}^{k}\left(\delta B_{t}-\frac{1}{2} \sum_{j=1}^{q} \sum_{n=0}^{\infty} \sum_{m=0}^{\infty}\left(\frac{\operatorname{ad}_{Y_{t}^{\varepsilon, \delta}}^{n}\left(\sigma_{t}^{j}\right)}{(n+1)!} \frac{\operatorname{ad}_{Y_{t}^{\varepsilon, \delta}}^{m}\left(\sigma_{t}^{j}\right)}{(m+1)!}+\frac{\left[\operatorname{ad}_{Y_{t}^{\varepsilon, \delta}}^{n}\left(\sigma_{t}^{j}\right), \operatorname{ad}_{Y_{t}^{\varepsilon, \delta}}^{m}\left(\sigma_{t}^{j}\right)\right]}{(n+m+2)(n+1)!m!}\right)\right)
\end{aligned}
$$

Equivalently, $Y^{\varepsilon, \delta}$ solves the Itô SDE

$$
\left\{\begin{array}{l}
d Y_{t}^{\varepsilon, \delta}=\mu^{\varepsilon, \delta}\left(t, Y_{t}^{\varepsilon, \delta}\right) d t+\sigma_{j}^{\varepsilon}\left(t, Y_{t}^{\varepsilon, \delta}\right) d W_{t}^{j}  \tag{2.11}\\
Y_{0}^{\varepsilon, \delta}=0
\end{array}\right.
$$

with

$$
\begin{aligned}
\sigma_{j}^{\varepsilon}(t, \cdot) & =\varepsilon \sum_{n=0}^{\infty} \frac{\beta_{n}}{n!} \operatorname{ad}^{n}\left(A_{t}^{(j)}\right), \quad j=1, \ldots, d, \\
\mu^{\varepsilon, \delta}(t, \cdot) & =\sum_{n=0}^{\infty} \frac{\beta_{n}}{n!} \operatorname{ad}^{n}\left(\delta B_{t}-\frac{1}{2} \sum_{j=1}^{q} Q .\left(\sigma_{j}^{\varepsilon}(t, \cdot), \sigma_{j}^{\varepsilon}(t, \cdot)\right)\right) .
\end{aligned}
$$

We now assume that $Y^{\varepsilon, \delta}$ admits the representation

$$
\begin{equation*}
Y_{t}^{\varepsilon, \delta}=\sum_{n=0}^{\infty} \sum_{r=0}^{n} Y_{t}^{(r, n-r)} \varepsilon^{r} \delta^{n-r} \tag{2.12}
\end{equation*}
$$

for a certain family $\left(Y^{(r, n-r)}\right)_{n, r \in \mathbb{N}_{0}}$ of stochastic processes. In particular, setting $(\varepsilon, \delta)=(1,1), 2.12$ would yield

$$
\begin{equation*}
Y_{t}=\sum_{n=0}^{\infty} Y_{t}^{(n)} \quad \text { with } \quad Y_{t}^{(n)}:=\sum_{r=0}^{n} Y_{t}^{(r, n-r)} \tag{2.13}
\end{equation*}
$$

Remark 2.3. Note that it is possible to re-order the double series $\sum_{n=0}^{\infty} \sum_{r=0}^{n} Y_{t}^{(r, n-r)}$ according to any arbitrary choice, for the latter will be proved to be absolutely convergent. The above choice for $Y^{(n)}$ contains all the terms of equal order by weighing $\varepsilon$ and $\delta$ in the same way. A different choice, which respects the probabilistic relation $\sqrt{\Delta t} \approx \Delta W_{t}$, corresponds to weighing $\delta$ as $\varepsilon^{2}$. This would lead to setting

$$
Y_{t}^{(n)}:=\sum_{r=0}^{\left\lfloor\frac{n}{2}\right\rfloor} Y_{t}^{(n-2 r, r)}
$$

in 2.13.
Remark 2.4. Observe that, if the function $(\varepsilon, \delta) \mapsto Y_{0}^{\varepsilon, \delta}$ is assumed to be continuous $P$-almost surely, than the initial condition in 2.11 implies

$$
Y_{0}^{(i, j)}=0 \quad P \text {-a.s., } \quad i, j \in \mathbb{N}_{0}
$$

and thus

$$
Y_{0}^{(n)}=0 \quad P \text {-a.s. }, \quad n \in \mathbb{N}_{0}
$$

We now plug 2.12 into 2.11 and collect all terms of equal order in $\varepsilon$ and $\delta$. Up to order 2 we obtain

$$
\begin{array}{rlrl}
\varepsilon^{0} \delta^{0}: & Y_{t}^{(0,0)} & =0 \\
\varepsilon^{1} \delta^{0}: & Y_{t}^{(1,0)} & =\int_{0}^{t} A_{s}^{(j)} d W_{s}^{j} \\
\varepsilon^{0} \delta^{1}: & Y_{t}^{(0,1)} & =\int_{0}^{t} B_{s} d s  \tag{2.16}\\
\varepsilon^{2} \delta^{0}: & Y_{t}^{(2,0)} & =-\frac{1}{2} \int_{0}^{t}\left(A_{s}^{(j)}\right)^{2} d s+\frac{1}{2} \int_{0}^{t}\left[A_{s}^{(j)}, \int_{0}^{s} A_{r}^{(i)} d W_{r}^{i}\right] d W_{s}^{j} \\
\varepsilon^{1} \delta^{1}: & Y_{t}^{(1,1)}=\frac{1}{2} \int_{0}^{t}\left[B_{s}, \int_{0}^{s} A_{r}^{(j)} d W_{r}\right] d s+\frac{1}{2} \int_{0}^{t}\left[A_{s}^{(j)}, \int_{0}^{s} B_{r} d r\right] d W_{s}^{j} \\
\varepsilon^{0} \delta^{2}: & Y_{t}^{(0,2)}=\frac{1}{2} \int_{0}^{t}\left[B_{s}, \int_{0}^{s} B_{r} d r\right] d s
\end{array}
$$

for any $t \geq 0$, where we used, one more time, Einstein summation convention to imply summation over the indexes $i, j$ and Remark 2.4 to set all the initial conditions equal to zero. Proceeding by induction, one can obtain a recursive representation for the general term $Y^{(r, n-r)}$ in 2.12), namely:

$$
\begin{equation*}
Y_{t}^{(r, n-r)}=\int_{0}^{t} \mu_{s}^{r, n-r} d s+\int_{0}^{t} \sigma_{s}^{r, n-r, j} d W_{s}^{j}, \quad n \in \mathbb{N}_{0}, r=0, \ldots, n \tag{2.17}
\end{equation*}
$$

where the terms $\sigma^{r, n-r, j}, \mu^{r, n-r}$ are defined recursively as

$$
\begin{align*}
\sigma_{s}^{r, n-r, j} & :=\sum_{i=0}^{n-1} \frac{\beta_{i}}{i!} S_{s}^{r-1, n-r, i}\left(A^{(j)}\right)  \tag{2.18}\\
\mu_{s}^{r, n-r} & :=\sum_{i=0}^{n-1} \frac{\beta_{i}}{i!} S_{s}^{r, n-r-1, i}(B)-\frac{1}{2} \sum_{j=1}^{q} \sum_{i=0}^{n-2} \frac{\beta_{i}}{i!} \sum_{q_{1}=2}^{r} \sum_{q_{2}=0}^{n-r} S^{r-q_{1}, n-r-q_{2}, i}\left(Q^{q_{1}, q_{2}, j}\right), \tag{2.19}
\end{align*}
$$

with

$$
\left.\left.\begin{array}{rl}
Q_{s}^{q_{1}, q_{2}, j}:=\sum_{i_{1}=2}^{q_{1}} \sum_{i_{2}=0}^{q_{2}} \sum_{h_{1}=1}^{i_{1}-1} \sum_{h_{2}=0}^{i_{2}} & \sum_{p_{1}=0}^{q_{1}-i_{1}} \sum_{p_{2}=0}^{q_{2}-i_{2}} \\
& \sum_{m_{1}=0}^{p_{1}+p_{2}} \frac{\sum_{m_{2}=0}^{q_{1}-i_{1}-p_{1}+q_{2}-i_{2}-p_{2}}}{\left(m_{1}+1\right)!}\left(\sigma_{s}^{p_{1}, p_{2}, m_{1}}\left(\sigma_{1}^{h_{1}, h_{2}, j}\right)\right. \\
S_{s}^{q_{1}-i_{1}-p_{1}, q_{2}-i_{2}-p_{2}, m_{2}}\left(\sigma_{s}^{i_{1}-h_{1}, i_{2}-h_{2}, j}\right) \\
\left(m_{2}+1\right)!
\end{array}\right), ~+\frac{\left[S_{s}^{p_{1}, p_{2}, m_{1}}\left(\sigma_{s}^{i_{1}-h_{1}, i_{2}-h_{2}, j}\right), S_{s}^{q_{1}-i_{1}-p_{1}, q_{2}-i_{2}-p_{2}, m_{2}}\left(\sigma_{s}^{h_{1}, h_{2}, j}\right)\right]}{\left(m_{1}+m_{2}+2\right)\left(m_{1}+1\right)!m_{2}!}\right)
$$

and with the operators $S$ being defined as

$$
\begin{aligned}
S_{s}^{r-1, n-r, 0}(A) & := \begin{cases}A & \text { if } r=n=1, \\
0 & \text { otherwise },\end{cases} \\
S_{s}^{r-1, n-r, i}(A) & :=\sum_{\substack{\left(j_{1}, k_{1}\right), \ldots,\left(j_{i}, k_{i}\right) \in \mathbb{N}_{0}^{2} \\
j_{1}+\cdots+j_{i}=r-1 \\
k_{1}+\cdots+k_{i}=n-r}}\left[Y_{s}^{\left(j_{1}, k_{1}\right)},\left[\ldots,\left[Y_{s}^{\left(j_{i}, k_{i}\right)}, A_{s}\right] \ldots\right]\right] \\
& =\sum_{\substack{\left(j_{1}, k_{1}\right), \ldots,\left(j_{i}, k_{i}\right) \in \mathbb{N}_{0}^{2} \\
j_{1}+\cdots+j_{i}=r-1 \\
k_{1}+\cdots k_{i}=n-r}} \operatorname{ad}_{Y_{s}^{\left(j_{1}, k_{1}\right)} \circ \cdots \circ \operatorname{ad}_{Y_{s}}^{\left(j_{i}, k_{i}\right)}\left(A_{s}\right), \quad i \in \mathbb{N} .}
\end{aligned}
$$

Remark 2.5. All the processes $Y^{(r, n-r)}$, with $n \in \mathbb{N}$ and $r=0, \ldots, n$, are well defined according to the recursion 2.17-2.18-2.19, as long as $B$ and $A^{(1)}, \ldots, A^{(q)}$ are bounded and progressively measurable stochastic processes.

Example 2.6. As we already pointed out in the introduction, in the case $j=1$ and $B \equiv 0$, the SDE (1.1) admits an explicit solution given by

$$
Y_{t}=-\frac{1}{2} A^{2} t+A W_{t}, \quad t \geq 0
$$

and the terms in the Magnus expansion 2.13 read as

$$
Y_{t}^{(1)}=A W_{t}, \quad Y_{t}^{(2)}=-\frac{1}{2} A^{2} t, \quad Y_{t}^{(n)}=0, n \geq 3
$$

In particular, the Magnus expansion coincides with the exact solution with the first two terms.

### 2.3 Examples and numerical tests

In this section we consider two examples, for which we compute the first three terms of the Magnus expansion given by 2.13 - 2.17 and present numerical experiments to test the accuracy of the approximate solutions to (1.1) stemming from it. In both cases we consider $j=1$ in (1.1) and replace $A^{(1)}$ with $A$ to shorten notation. The first example will be for constant matrices $A$ and $B$. In the second one we will consider $B \equiv 0$ and a deterministic upper diagonal $A_{t}$. For each numerical test we will implement the exponential of the truncated Magnus expansion up to order $n=1,2$ and 3 , i.e.

$$
\begin{equation*}
X^{(n)}:=e^{\sum_{i=1}^{n} Y^{(i)}}, \quad n=1,2,3 \tag{2.20}
\end{equation*}
$$

and compare it with a benchmark solution to (1.1).

Error and notations. Throughout this section we will employ the following tags:

- euler for the solution obtained with Euler-Maruyama scheme;
- exact to denote the time-discretization of an explicit solution, if available;
- m1, m2 and m3 for the time-discretization of the Magnus approximations in 2.20, up to order 1,2 and 3 , respectively.

For the error analysis we will make use of the following norms. Denoting by $X^{\text {ref }}$ and by $X^{\text {app }}$ a benchmark and an approximate solution, respectively, to 1.1 ) and by $\left(t_{k}\right)_{k=0, \ldots, N}$ a homogeneous discretization of $[0, t]$, we consider the random variable

$$
\operatorname{Err}_{t}:=\frac{\Delta}{t} \sum_{k=1}^{N} \frac{\left|X_{t_{k}}^{\mathrm{ref}}-X_{t_{k}}^{\mathrm{app}}\right|}{\left|X_{t_{k}}^{\mathrm{ref}}\right|} \approx \frac{1}{t} \int_{0}^{t} \frac{\left|X_{s}^{\mathrm{ref}}-X_{s}^{\mathrm{app}}\right|}{\left|X_{s}^{\mathrm{ref}}\right|} d s \quad \text { with } \Delta=\frac{t}{N}
$$

namely a discretization of the time-averaged relative error on the interval $[0, t]$. This is a way to measure the error on the whole trajectory as opposed to the error at a specific given time. Then we use Monte Carlo simulation, with $M$ independent realizations of the discretized Brownian trajectories, to approximate the distribution of $\operatorname{Err}_{t}$.

The matrix norm above is the Frobenius norm. In the following tests, $\mathrm{m} 1, \mathrm{~m} 2$ and m 3 will always play the role of $X^{\text {app }}$, exact always the role of $X^{\text {ref }}$, whereas euler will be either $X^{\text {app }}$ or $X^{\text {ref }}$ depending on whether exact is available or not.

We used for the calculations Matlab R2019a with Parallel Computing Toolbox running on Windows 10 Pro, on a machine with the following specifications: processor Intel(R) Core(TM) i7-8750H @ 2.20 GHz , $2 x 16$ GB (Dual Channel) Samsung SODIMM RAM @ 2667 MHz , and a NVIDIA GeForce RTX 2070 with Max-Q Design (8 GB GDDR6 RAM).

Also, we will make use of the Matlab built-in routine expm for the computation of the matrix exponential. As it turns out, this represents the most expensive step in the implementation of the Magnus approximation. However important, the pursue of optimized method for the matrix exponentiation is an extended topic of separate interest, which goes beyond the goals of this paper. Therefore, here we will limit ourselves to pointing out, separately, the computational times for the approximations of the logarithm and of the matrix exponential.

| Method | $t=0.25$ | $t=0.5$ | $t=0.75$ | $t=1$ | $t=2$ | $t=3$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| m 1 | $4.36 \%$ | $9.47 \%$ | $14.8 \%$ | $20.8 \%$ | $48.3 \%$ | $88.4 \%$ |
| m 2 | $0.177 \%$ | $0.471 \%$ | $0.888 \%$ | $1.55 \%$ | $5.2 \%$ | $10.7 \%$ |
| m 3 | $0.125 \%$ | $0.205 \%$ | $0.312 \%$ | $0.486 \%$ | $1.95 \%$ | $5.08 \%$ |

Table 1: $A$ and $B$ constant. Values of $E\left[\operatorname{Err}_{t}\right]$ (in percentage) for $\mathrm{m} 1, \mathrm{~m} 2$, m 3 , with euler as benchmark solution, obtained with $10^{3}$ samples.

Example: constant $A$ and $B$. With a slight abuse of notation, we consider $A_{t} \equiv A$ and $B_{t} \equiv B$. Recall that, if $A$ and $B$ do not commute, there is in general no closed-form solution to 1.1 . The first three terms of the Magnus expansion read as

$$
\begin{align*}
Y_{t}^{(1)}= & B t+A W_{t}, \quad Y_{t}^{(2)}=[A, B]\left(\frac{1}{2} t W_{t}-\int_{0}^{t} W_{s} d s\right)-\frac{1}{2} A^{2} t \\
Y_{t}^{(3)}= & {[[B, A], A]\left(\frac{1}{2} \int_{0}^{t} W_{s}^{2} d s-\frac{1}{2} W_{t} \int_{0}^{t} W_{s} d s+\frac{1}{12} t W_{t}^{2}\right) } \\
& +[[B, A], B]\left(\int_{0}^{t} s W_{s} d s-\frac{1}{2} t \int_{0}^{t} W_{s} d s-\frac{1}{12} t^{2} W_{t}\right) \tag{2.21}
\end{align*}
$$

We point out that, in this case, all the stochastic integrals appearing in the Magnus expansion can be solved in terms of Lebesgue integrals. Therefore, in order to discretize $Y^{(n)}$ it is not necessary to approximate stochastic integrals. This allows to use a sparser time grid compared to the Euler method, for which the discretization of stochastic integrals is necessary. In the following numerical tests, we discretize in time with mesh $\Delta$ equal to $10^{-4}$ for euler and equal to $10^{-3}$ for $\mathrm{m} 1, \mathrm{~m} 2$ and m 3 . Note that, as it is confirmed by the results in Table 1, choosing a finer time-discretization for euler (our reference method here) is essential in order to make it comparable with m3. Also, we choose $A$ and $B$ at random and normalize them by their spectral norms. In particular, the results below refer to

$$
A=\left(\begin{array}{cc}
0.335302 & -0.645492 \\
-0.264419 & 0.634641
\end{array}\right), \quad B=\left(\begin{array}{cc}
-0.0572262 & 0.0493763 \\
-0.665366 & 0.742744
\end{array}\right)
$$

In Figure 1 we plot one realization of the trajectories of the top-left component $\left(X_{t}\right)_{11}$, computed with the methods above, up to time $t=0.75$. In Table 1 we show the expectations $E\left[\operatorname{Err}_{t}\right]$ for different values of $t$, with euler as benchmark solution, computed via Monte Carlo simulation with $10^{3}$ samples. The same samples are used in Figure 2.3 to plot the empirical CDF of $\operatorname{Err}_{t}$.

The computational time for $10^{3}$ sampled trajectories, up to time $t=1$, is roughly 84 seconds for euler and 48 seconds for either $\mathrm{m} 1, \mathrm{~m} 2$ or m 3 . The latter, however, is divided as follows: nearly 2 seconds to compute the Magnus expansion and nearly 46 seconds to compute the matrix exponential with the Matlab function expm. We stress the fact that the computation of the logarithm via Magnus expansion is very fast because the computation of the integrals in $Y^{(1)}, Y^{(2)}$ and $Y^{(3)}$ can be parallelized.

Example: $j=1, B=0$ and $A_{t}$ upper triangular. We set

$$
A_{t}=\left(\begin{array}{cc}
2 & t  \tag{2.22}\\
0 & -1
\end{array}\right), \quad B_{t} \equiv 0
$$



Figure 1: $A$ and $B$ constant. One realization of the trajectories of the top-left component $\left(X_{t}\right)_{11}$, computed with euler, m1, m2, m3.


Figure 2: $A$ and $B$ constant. Empirical CDF of $\operatorname{Err}_{t}$, at $t=0.75$, for $\mathrm{m} 1, \mathrm{~m} 2$, m 3 , with euler as benchmark solution, obtained with $10^{3}$ samples.

| Method | $t=0.25$ | $t=0.5$ | $t=0.75$ | $t=1$ | $t=2$ | $t=3$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| euler | $0.484 \%$ | $0.708 \%$ | $0.884 \%$ | $1.07 \%$ | $2.07 \%$ | $3.15 \%$ |
| m1 | $18.9 \%$ | $41.9 \%$ | $69.2 \%$ | $105 \%$ | $405 \%$ | $1490 \%$ |
| m2 | $0.146 \%$ | $0.676 \%$ | $1.78 \%$ | $3.15 \%$ | $10.8 \%$ | $21.2 \%$ |
| m3 | $0.0576 \%$ | $0.252 \%$ | $0.759 \%$ | $1.44 \%$ | $6.01 \%$ | $14.1 \%$ |

Table 2: $B_{t}$ and $A_{t}$ as in 2.22. Values of $E\left[\operatorname{Err}_{t}\right]$ (in percentage) for euler, $\mathrm{m} 1, \mathrm{~m} 2, \mathrm{~m} 3$, with exact as benchmark solution, obtained with $10^{3}$ samples.

In this case 1.1 admits an explicit solution, which can be obtained by using Yoeurp and Yor's formula, given by

$$
X_{t}=\left(\begin{array}{cc}
e^{2\left(W_{t}-t\right)} & e^{2\left(W_{t}-t\right)}\left(\int_{0}^{t} s e^{-3 W_{s}+\frac{3}{2} s} d W_{s}-2 \int_{0}^{t} s e^{-3 W_{s}+\frac{3}{2} s} d s\right) \\
0 & e^{-\left(W_{t}+\frac{1}{2} t\right)}
\end{array}\right)
$$

The first three terms of the Magnus expansion read as

$$
\begin{aligned}
& Y_{t}^{(1)}=\left(\begin{array}{cc}
2 W_{t} & t W_{t}-\int_{0}^{t} W_{s} d s \\
0 & -W_{t}
\end{array}\right), \quad Y_{t}^{(2)}=-\frac{1}{2}\left(\begin{array}{cc}
4 t & \frac{1}{2} t^{2} \\
0 & t
\end{array}\right)-\frac{3}{2}\left(\begin{array}{cc}
0 & W_{t} \int_{0}^{t} W_{u} d u-\int_{0}^{t} W_{u}^{2} d u \\
0 & 0
\end{array}\right), \\
& Y_{t}^{(3)}=\left(\begin{array}{cc}
0 & \frac{3}{4}\left(t-W_{t}^{2}\right) \int_{0}^{t} W_{u} d u-\frac{3}{2} \int_{0}^{t} W_{u} u d u+\frac{9}{4} W_{t} \int_{0}^{t} W_{u}^{2} d u-\frac{3}{2} \int_{0}^{t} W_{u}^{3} d u+\frac{3}{8} t^{2} W_{t} \\
0 & 0
\end{array}\right) .
\end{aligned}
$$

Again, all the stochastic integrals appearing in the Magnus expansion can be solved in terms of Lebesgue integrals, which allows us to use one more time a sparser time grid compared to the Euler method and the discretized exact solution. In the following numerical tests, we discretize in time with mesh $\Delta$ equal to $10^{-4}$ for exact and euler, and equal to $10^{-3}$ for $\mathrm{m} 1, \mathrm{~m} 2$ and m 3 . Note that euler serves here as an alternative approximation and that choosing a finer time-discretization for euler and exact (our reference method here) is again essential in order to make them comparable with m 3 .

In Table 2 we show the expectations $E\left[\operatorname{Err}_{t}\right]$ for different values of $t$, with exact as benchmark solution, computed via Monte Carlo simulation with $10^{3}$ samples. The same samples are used in Figure 3 to plot the empirical CDF of $\operatorname{Err}_{t}$.

The computational time for $10^{3}$ sampled trajectories, up to time $t=1$, is approximately 9.1 seconds for exact, 48 seconds for euler and 40 seconds for either $\mathrm{m} 1, \mathrm{~m} 2$ or m 3 . The latter, however, is divided as follows: nearly 0.1 seconds to compute the Magnus expansion and nearly 40 seconds to compute the matrix exponential with the Matlab function expm. Again, note that the computation of the logarithm via Magnus expansion is very fast thanks to the possibility of parallelizing the computation of the integrals in $Y^{(1)}, Y^{(2)}$ and $Y^{(3)}$.

As it appears in the results above, the accuracy of the Magnus expansion quickly deteriorates as the time increases. This is largely due the fact that the spectral norm

$$
\left\|A_{t}\right\|=\sqrt{\frac{1}{2}\left(t^{2}+\sqrt{t^{4}+10 t^{2}+9}+5\right)}
$$

is an increasing function of $t$. This behavior shall not come as a surprise, since the proof of Theorem 1.1 already uncovered the relation between the convergence time $\tau$ and the spectral norms of $A_{t}$ and $B_{t}$. Such


Figure 3: $B_{t}$ and $A_{t}$ as in 2.22. Empirical CDF of $\operatorname{Err}_{t}$, at $t=0.75$, for euler, $\mathrm{m} 1, \mathrm{~m} 2$, m3, with exact as benchmark solution, obtained with $10^{3}$ samples.

| Method | $t=0.25$ | $t=0.5$ | $t=0.75$ | $t=1$ | $t=2$ | $t=3$ | $t=10$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| euler | $0.128 \%$ | $0.183 \%$ | $0.226 \%$ | $0.256 \%$ | $0.383 \%$ | $0.471 \%$ | $0.682 \%$ |
| m1 | $4.46 \%$ | $9.03 \%$ | $13.4 \%$ | $18.2 \%$ | $33.6 \%$ | $47.6 \%$ | $94 \%$ |
| m2 | $0.0214 \%$ | $0.102 \%$ | $0.257 \%$ | $0.478 \%$ | $2.08 \%$ | $3.89 \%$ | $15.2 \%$ |
| m3 | $0.0118 \%$ | $0.035 \%$ | $0.0756 \%$ | $0.134 \%$ | $0.621 \%$ | $1.21 \%$ | $5.41 \%$ |

Table 3: $B \equiv 0$ and $A_{t}$ as in 2.22 normalized by its spectral norm. Values of $E\left[\operatorname{Err}_{t}\right]$ (in percentage) for euler, m1, m2, m3, with exact as benchmark solution, obtained with $10^{3}$ samples.
relation is also consistent with the convergence condition 1.7) that holds in the deterministic case. In order to asses numerically the impact of the spectral norm of $A_{t}$ on the quality of the Magnus approximation, we now repeat the experiments on the equation obtained by normalizing $A_{t}$ as in 2.22 with respect to $\left\|A_{t}\right\|$. As it turns out, the accuracy of $\mathrm{m} 1, \mathrm{~m} 2$ and m 3 improves considerably with this normalization. Note that, in this case, 1.1 no longer admits a closed-form solution, while the representation for the terms $Y^{(1)}, Y^{(2)}$ and $Y^{(3)}$ in the Magnus expansion is omitted for it becomes rather tedious to write. In Figure 4 we plot one realization of the trajectories of the top-right component $\left(X_{t}\right)_{12}$, computed with all the methods above, up to time $t=10$. In this case we did not plot a diagonal component of the solution because the latter are exact for m 2 and m 3 , up to discretization errors of Lebesgue integrals. Table 3 and Figure 5 are analogous to Table 2 and Figure 3 and are obtained again with $10^{3}$ independent samples. Computational times are now roughly 2 seconds higher for each method, due to the increased complexity of the computations.


Figure 4: $B \equiv 0$ and $A_{t}$ as in 2.22 normalized by its spectral norm. One realization of the trajectories of the top-right component $\left(X_{t}\right)_{12}$, computed with exact, euler, m1, m2, m3.


Figure 5: $B \equiv 0$ and $A_{t}$ as in 2.22 normalized by its spectral norm. Empirical CDF of $\operatorname{Err}_{t}$, at $t=0.75$, for euler, m1, m2, m3, with exact as benchmark solution, obtained with $10^{3}$ samples.

## 3 Convergence analysis

In this section we prove Theorem 1.1. To avoid ambiguity, only in this section, we denote by $\mathcal{M}_{\mathbb{R}}^{d \times d}$ and $\mathcal{M}_{\mathbb{C}}^{d \times d}$ the spaces of $(d \times d)$-matrices with real and complex entries, respectively; on these spaces we shall make use of the Frobenius norm denoted by $|\cdot|$. We say that a matrix-valued function is holomorphic if all its entries are holomorphic functions. We recall that $W=\left(W^{1}, \ldots, W^{q}\right)$ is a $q$-dimensional standard Brownian motion and $A^{(1)}, \ldots, A^{(q)}, B$ are bounded $\mathcal{M}_{\mathbb{R}}^{d \times d}$-valued progressively measurable stochastic processes defined on a filtered probability space $\left(\Omega, \mathcal{F}, P,\left(\mathcal{F}_{t}\right)_{t \geq 0}\right)$. Also recall that, for any $\mathcal{M}_{\mathbb{R}}^{d \times d}$-valued process $M=\left(M_{t}\right)_{t \in[0, T]}$, we set $\|M\|_{T}:=\||M|\|_{L^{\infty}([0, T] \times \Omega)}$.

We start with two preliminary lemmas.
Lemma 3.1. Assume that $Y=\left(Y_{t}^{\varepsilon, \delta}\right)_{\varepsilon, \delta \in \mathbb{R}, t \in \mathbb{R}_{\geq 0}}$ is a $\mathcal{M}_{\mathbb{R}}^{d \times d}$-valued stochastic process of the form 2.12. If $Y$ solves the $S D E$ 2.11 up to a positive stopping time $\tau$, then $Y^{(r, n-r)}$ in 2.12 are the Ito processes defined through 2.17-2.18-2.19 for any $t<\tau$.

Proof. We prove 2.17)-2.18-2.19 only for $n=0,1$. Namely, we show that 2.14, 2.15 and 2.16 hold up to time $\tau, P$-a.s. The representation for the general term $Y^{(r, n-r)}$ can be proved by induction; we omit the details for brevity.

Since $Y$ is of the form (2.12) then $Y_{t}^{(0,0)}=Y_{t}^{0,0}$ for any $t<\tau$. Moreover, since $Y$ solves the SDE 2.11) then $Y^{0,0} \equiv 0$ on $[0, \tau[, P$-a.s. Thus 2.14$]$ holds up to time $\tau, P$-a.s.

Now, (2.11) yields

$$
\begin{equation*}
Y_{t}^{\varepsilon, 0}=\varepsilon \int_{0}^{t} A_{s}^{(j)} d W_{s}^{j}+\varepsilon R_{t}^{\varepsilon}, \quad t \in[0, \tau[, \quad P \text {-a.s. } \tag{3.1}
\end{equation*}
$$

where

$$
R_{t}^{\varepsilon}=\int_{0}^{t}\left(\sum_{k=1}^{\infty} \frac{\beta_{k}}{k!} \operatorname{ad}_{Y_{s}^{\varepsilon, 0}}^{k}\left(A_{s}^{(j)}\right)\right) d W_{s}^{j}-\frac{\varepsilon}{2} \int_{0}^{t} \mathcal{L}_{Y_{s}^{\varepsilon, 0}}^{-1}\left(Q_{Y_{s}^{\varepsilon, 0}}\left(\sum_{k=0}^{\infty} \frac{\beta_{k}}{k!} \operatorname{ad}_{Y_{s}^{\varepsilon, 0}}^{k}\left(A_{s}^{(j)}\right), \sum_{k=0}^{\infty} \frac{\beta_{k}}{k!} \operatorname{ad}_{Y_{s}^{\varepsilon, 0}}^{k}\left(A_{s}^{(j)}\right)\right)\right) d s
$$

Note that, again by 2.11, $R^{0} \equiv 0$-a.s. Moreover, representation 2.12 implies continuity of $\varepsilon \mapsto Y_{t}^{\varepsilon, 0}$ near $\varepsilon=0$, which in turn implies the continuity of $\varepsilon \mapsto R_{t}^{\varepsilon}$. Thus we have $\lim _{\varepsilon \rightarrow 0} R_{t}^{\varepsilon}=R_{t}^{0} P$-a.s. This, together with (3.1) and 2.12 implies that 2.15 necessarily holds, up to time $\tau, \stackrel{\varepsilon \rightarrow 0}{P \rightarrow \text {-a.s. }}$

Similarly, 2.11 yields

$$
Y_{t}^{0, \delta}=\delta \int_{0}^{t} B_{s} d s+\delta Q_{t}^{\delta}, \quad t \in[0, \tau[, \quad P \text {-a.s. }
$$

with

$$
Q_{t}^{\delta}=\int_{0}^{t}\left(\sum_{k=1}^{\infty} \frac{\beta_{k}}{k!} \mathrm{ad}_{Y_{s}^{0, \delta}}^{k}\left(B_{s}\right)\right) d s
$$

and the same argument employed above yields 2.16 up to time $\tau, P$-almost surely.
Lemma 3.2. Let $M \in \mathcal{M}_{\mathbb{C}}^{d \times d}$ be nonsingular and such that $\left\|M-I_{d}\right\|<1$ where $\|\cdot\|$ is the spectral norm. Then $M$ has a unique logarithm, which is

$$
\log M=\sum_{n=1}^{\infty}(-1)^{n+1} \frac{\left(M-I_{d}\right)^{n}}{n}
$$

$$
=\left(M-I_{d}\right) \int_{0}^{\infty} \frac{1}{1+\mu}\left(\mu I_{d}+M\right)^{-1} d \mu
$$

In particular, we have

$$
\begin{equation*}
\|\log M\| \leq-\log \left(1-\left\|M-I_{d}\right\|\right) \tag{3.2}
\end{equation*}
$$

Proof. The first representation is a standard result. The second representation stems from the factorization $M=V J V^{-1}$ with $J$ in Jordan form, under the assumption that $M$ has no non-positive real eigenvalues, i.e. $\lambda \in \mathbb{C} \backslash]-\infty, 0]$ for any $\lambda$ eigenvalue of $M$. This last property, however, is ensured by the assumption $\left\|M-I_{d}\right\|<1$. Indeed, the latter implies

$$
|M v-v|<1, \quad v \in \mathbb{R}^{d},|v|=1
$$

which in turn implies that, if $\lambda$ is a real eigenvalue of $M$ and $v$ is one of its normalized eigenvectors, then

$$
1>|M v-v|=|\lambda v-v|=|\lambda-1| .
$$

We have one last preliminary lemma, containing some technical results concerning the solutions to 2.7. These are semi-standard, in that they can be inferred by combining and adapting existing results in the literature.
Lemma 3.3. For any $T>0$ and $\varepsilon, \delta \in \mathbb{C}$, the $S D E$ 2.7) has a unique strong solution $\left(X_{t}^{\varepsilon, \delta}\right)_{t \in[0, T]}$. For any $p \geq 1$ and $h>0$ there exists a positive constant $\kappa$, only dependent on $\left\|A^{(1)}\right\|_{T}, \ldots,\left\|A^{(q)}\right\|_{T},\|B\|_{T}$, $d$, $T, h$ and $p$, such that

$$
\begin{align*}
E\left[\left|X_{t}^{\varepsilon, \delta}-X_{s}^{\varepsilon^{\prime}, \delta^{\prime}}\right|^{2 p}\right] & \leq \kappa\left(|t-s|^{p}+\left(\left|\varepsilon-\varepsilon^{\prime}\right|+\left|\delta-\delta^{\prime}\right|\right)^{2 p}\right)  \tag{3.3}\\
E\left[\sup _{0 \leq u \leq t}\left|X_{u}^{\varepsilon, \delta}-X_{0}^{\varepsilon, \delta}\right|^{2 p}\right] & \leq \kappa t^{p}(|\varepsilon|+|\delta|)^{2 p} \tag{3.4}
\end{align*}
$$

for any $0 \leq t, s \leq T$ and $\varepsilon, \delta, \varepsilon^{\prime}, \delta^{\prime} \in \mathbb{C}$ with $|\varepsilon|,|\delta|,\left|\varepsilon^{\prime}\right|,\left|\delta^{\prime}\right| \leq h$.
Up to modifications, $\left(X_{t}^{\varepsilon, \delta}\right)_{\varepsilon, \delta \in \mathbb{C}, t \in[0, T]}$ is a continuous process such that:
i) for any $t \in[0, T]$, the function $(\varepsilon, \delta) \mapsto X_{t}^{\varepsilon, \delta}$ is holomorphic;
ii) the functions $(t, \varepsilon, \delta) \mapsto \partial_{\varepsilon} X_{t}^{\varepsilon, \delta}$ and $(t, \varepsilon, \delta) \mapsto \partial_{\delta} X_{t}^{\varepsilon, \delta}$ are continuous;
iii) for any $p \geq 1$ and $h>0$ there exists a positive constant $\kappa$ only dependent on $\left\|A^{(1)}\right\|_{T}, \ldots,\left\|A^{(q)}\right\|_{T}$, $\|B\|_{T}, d, T, h$ and $p$, such that

$$
\begin{equation*}
E\left[\sup _{0 \leq s \leq t}\left\{\left|\partial_{\varepsilon} X_{s}^{\varepsilon, \delta}\right|^{2 p}+\left|\partial_{\delta} X_{s}^{\varepsilon, \delta}\right|^{2 p}\right\}\right] \leq \kappa t^{p}(|\varepsilon|+|\delta|)^{p} \tag{3.5}
\end{equation*}
$$

for any $t \in[0, T]$ and $|\varepsilon|,|\delta| \leq h$.
Proof. Existence of the solution and estimates (3.3)-3.4) of the moments follow from the results in Section 5 , Chapter 2 in [14] (in particular, see Corollary 5 on page 80 and Theorem 7 on page 82).

The second part of the statement is a refined version of the Kolmogorov continuity theorem in the form that can be found for instance in Section 2.3 in [17]: a detailed proof is provided in [12].

Remark 3.4. The existence and uniqueness for the solution to 1.1 is a particular case of the previous result.

We are now in the position to prove Theorem 1.1.
Proof of Theorem 1.1. We fix $h>1, T>0$, and let $\left(X_{t}^{\varepsilon, \delta}\right)_{\varepsilon, \delta \in \mathbb{C}, t \in[0, T]}$ be the solution of the SDE (2.7) as defined in Lemma 3.3. Moreover, for $t \in] 0, T]$, we set $\left.Q_{t, h}:=\right] 0, t\left[\times B_{h}(0)\right.$ where $B_{h}(0)=\left\{(\varepsilon, \delta) \in \mathbb{C}^{2} \mid\right.$ $|(\varepsilon, \delta)|<h\}$.


$$
\begin{equation*}
\tau:=\sup \left\{t \in[0, T]| | X_{s}^{\varepsilon, \delta}-I_{d} \mid<1-e^{-\pi} \text { for any }(s, \varepsilon, \delta) \in Q_{t, h}\right\} \tag{3.6}
\end{equation*}
$$

is strictly positive. Furthermore, again by continuity,

$$
(\tau \leq t)=\bigcup_{(s, \varepsilon, \delta) \in \widetilde{Q}_{t, h}}\left(\left|X_{s}^{\varepsilon, \delta}-I_{d}\right| \geq 1-e^{-\pi}\right), \quad t \in[0, T]
$$

where $\widetilde{Q}_{t, h}$ is a countable, dense subset of $Q_{t, h}$, which implies that $\tau$ is a stopping time.
Let $(t, \varepsilon, \delta) \in Q_{\tau, h}$ : by Lemma 3.2 applied to $M=X_{t}^{\varepsilon, \delta}$ we have

$$
\begin{equation*}
Y_{t}^{\varepsilon, \delta}:=\log X_{t}^{\varepsilon, \delta}=\sum_{n=1}^{\infty}(-1)^{n+1} \frac{\left(X_{t}^{\varepsilon, \delta}-I_{d}\right)^{n}}{n}=\left(X_{t}^{\varepsilon, \delta}-I_{d}\right) \int_{0}^{\infty} \frac{1}{1+\mu}\left(\mu I_{d}+X_{t}^{\varepsilon, \delta}\right)^{-1} d \mu \tag{3.7}
\end{equation*}
$$

Notice that $X_{t}^{\varepsilon, \delta}$ (and therefore also $Y_{t}^{\varepsilon, \delta}$ ) is real for $\varepsilon, \delta \in \mathbb{R}$ : in particular, $Y_{t}=Y_{t}^{1,1}$ is real and this proves Part (i).

Part (ii): since $(\varepsilon, \delta) \mapsto X_{t}^{\varepsilon, \delta}$ is holomorphic, we can differentiate 3.7 to infer that $(\varepsilon, \delta) \mapsto Y_{t}^{\varepsilon, \delta}$ is holomorphic as well: indeed, we have

$$
\begin{aligned}
\partial_{\varepsilon} Y_{t}^{\varepsilon, \delta}= & \partial_{\varepsilon} X_{t}^{\varepsilon, \delta} \int_{0}^{\infty} \frac{1}{1+\mu}\left(\mu I_{d}+X_{t}^{\varepsilon, \delta}\right)^{-1} d \mu \\
& +\left(X_{t}^{\varepsilon, \delta}-I_{d}\right) \int_{0}^{\infty} \frac{1}{1+\mu}\left(\mu I_{d}+X_{t}^{\varepsilon, \delta}\right)^{-1}\left(\partial_{\varepsilon} X_{t}^{\varepsilon, \delta}\right)\left(\mu I_{d}+X_{t}^{\varepsilon, \delta}\right)^{-1} d \mu, \quad(t, \varepsilon, \delta) \in Q_{\tau, h}
\end{aligned}
$$

and similarly by differentiating w.r.t. to $\delta$. Then the expansion of $Y_{t}^{\varepsilon, \delta}$ in power series at $(\varepsilon, \delta)=(0,0)$ is absolutely convergent on $B_{h}(0)$ and the representation 2.12 holds on $Q_{\tau, h}$ for some random coefficients $Y_{t}^{(r, n-r)}$. To conclude we need to show that the latter are as given by 2.17)-2.18- 2.19). Then (1.9) will stem from 2.12 by setting $(\varepsilon, \delta)=(1,1)$. In light of Lemma 3.2 , the logarithmic map is continuously twice differentiable on the open subset of $\mathcal{M}_{\mathbb{C}}^{d \times d}$ of the matrices $M$ such that $\left\|M-I_{d}\right\|<1$ : thus $Y_{t}^{\varepsilon, \delta}$ admits an Itô representation (2.5) for $(t, \varepsilon, \delta) \in Q_{\tau, h}$. Then Proposition 2.1 together with 2.7 yield 2.9-2.10 P-a.s. up to $\tau$ for any $(\varepsilon, \delta) \in B_{h}(0) \cap \mathbb{R}^{2}$. Furthermore, by estimate 3.2 of Lemma 3.2 we also have $\left\|Y_{t}^{\varepsilon, \delta}\right\|<\pi$ for $t<\tau$. Therefore, we can apply Baker's Lemma 2.2 to invert $\mathcal{L}_{Y_{t}^{\varepsilon, \delta}}$ in 2.9- 2.10 and obtain that $Y^{\varepsilon, \delta}$ solves 2.11 up to $\tau$ for any $(\varepsilon, \delta) \in B_{h}(0) \cap \mathbb{R}^{2}$. Part (ii) then follows from Lemma 3.1.

Part (iii): for $t \leq T$ let

$$
f_{t}(\varepsilon, \delta):=\max _{s \in[0, t]}\left|X_{s}^{\varepsilon, \delta}-I_{d}\right|, \quad M_{t}:=\sup _{(\varepsilon, \delta) \in B_{h}(0)} f_{t}(\varepsilon, \delta) .
$$

By definition 3.6, we have

$$
\begin{equation*}
P(\tau \leq t) \leq P\left(M_{t} \geq 1-e^{-\pi}\right) \leq \frac{1}{\left(1-e^{-\pi}\right)^{2}} E\left[M_{t}^{2}\right] \tag{3.8}
\end{equation*}
$$

and therefore 1.10 follows by suitably estimating $E\left[M_{t}^{2}\right]$. To prove such an estimate we will show in the last part of the proof that $f_{t}$ belongs to the Sobolev space $W^{1,2 p}\left(B_{h}(0)\right)$ for any $p \geq 1$ and we have

$$
\begin{equation*}
E\left[\left\|f_{t}\right\|_{W^{1,2 p}\left(B_{h}(0)\right)}^{2 p}\right] \leq C t^{p}, \quad t \in[0, T] \tag{3.9}
\end{equation*}
$$

where the positive constant $C$ depends only on $\left\|A^{(1)}\right\|_{T}, \ldots,\left\|A^{(q)}\right\|_{T},\|B\|_{T}, d, T, h$ and $p$. Since $f_{t} \in$ $W^{1,2 p}\left(B_{h}(0)\right)$ and $B_{h}(0) \subseteq \mathbb{R}^{4}$, by Morrey's inequality (cf., for instance, Corollary 9.14 in [4]) for any $p>2$ we have

$$
\begin{equation*}
M_{t} \leq c_{0}\left\|f_{t}\right\|_{W^{1,2 p}\left(B_{h}(0)\right)} \tag{3.10}
\end{equation*}
$$

where $c_{0}$ is a a positive constant, dependent only on $p$ and $h$ (in particular, $c_{0}$ is independent of $\omega$ ). Combining (3.9) with 3.10, for a fixed $p>2$ we have

$$
E\left[M_{t}^{2}\right] \leq c_{0}^{2} E\left[\left\|f_{t}\right\|_{W^{1,2 p}\left(B_{h}(0)\right)}^{2}\right] \leq
$$

(by Hölder inequality)

$$
\leq c_{0}^{2} C t, \quad t \in[0, T]
$$

This last estimate, combined with (3.8), proves (1.10).
To conclude, we are left with the proof of 3.9 . First we have

$$
\begin{equation*}
E\left[\int_{B_{h}(0)}\left|f_{t}(\varepsilon, \delta)\right|^{2 p} d \varepsilon d \delta\right]=\int_{B_{h}(0)} E\left[\left|f_{t}(\varepsilon, \delta)\right|^{2 p}\right] d \varepsilon d \delta \leq C t^{p} \tag{3.11}
\end{equation*}
$$

where we used the estimate (3.4) of Lemma 3.3 in the last inequality. Fix now $t \in] 0, T],(\varepsilon, \delta),\left(\varepsilon^{\prime}, \delta^{\prime}\right) \in B_{h}(0)$ such that $f_{t}\left(\varepsilon^{\prime}, \delta^{\prime}\right) \leq f_{t}(\varepsilon, \delta)$ and set

$$
\bar{t}=\underset{0 \leq s \leq t}{\arg \max }\left|X_{s}^{\varepsilon, \delta}-I_{d}\right|, \quad \tilde{t}=\underset{0 \leq s \leq t}{\arg \max }\left|X_{s}^{\varepsilon^{\prime}, \delta^{\prime}}-I_{d}\right| .
$$

Note that the $\arg$ max above do exist in that the process $g_{s}(\varepsilon, \delta):=X_{s}^{\varepsilon, \delta}-I_{d}$ is continuous in $s$ and we have
where $\nabla=\nabla_{\varepsilon, \delta}$. This, as $(s, \varepsilon, \delta) \mapsto \nabla g_{s}(\varepsilon, \delta)$ is continuous on $Q_{t, h}$, implies $f_{t} \in W^{1,2 p}\left(B_{h}(0)\right)$ and yields the key inequality

$$
\left|\nabla f_{t}(\varepsilon, \delta)\right| \leq \sup _{0 \leq s \leq t}\left|\nabla X_{s}^{\varepsilon, \delta}\right|, \quad(\varepsilon, \delta) \in B_{h}(0)
$$

Therefore, we have

$$
E\left[\int_{B_{h}(0)}\left|\nabla f_{t}(\varepsilon, \delta)\right|^{2 p} d \varepsilon d \delta\right]=\int_{B_{h}(0)} E\left[\left|\nabla f_{t}(\varepsilon, \delta)\right|^{2 p}\right] d \varepsilon d \delta \leq \int_{B_{h}(0)} E\left[\sup _{0 \leq s \leq t}\left|\nabla X_{s}^{\varepsilon, \delta}\right|^{2 p}\right] d \varepsilon d \delta \leq C t^{p},
$$

where we used the estimate (3.5) of Lemma 3.3 in the last inequality. This, together with 3.11), proves (3.9) and conclude the proof.

## 4 Applications to SPDEs

The aim of this section is to apply the stochastic Magnus expansion for the numerical solution of parabolic stochastic partial differential equations (SPDEs).

### 4.1 Stochastic Cauchy problem and fundamental solution

Let $\left(\Omega, \mathcal{F}, P,\left(\mathcal{F}_{t}\right)_{t \geq 0}\right)$ be a filtered probability space endowed with a real Brownian motion $W$. We consider the stochastic Cauchy problem

$$
\left\{\begin{array}{l}
d u_{t}(x)=\mathbf{L}_{t} u_{t}(x) d t+\mathbf{G}_{t} u_{t}(x) d W_{t}, \quad t>0, x \in \mathbb{R}  \tag{4.1}\\
u_{0}=\varphi
\end{array}\right.
$$

where $\mathbf{L}_{t}$ is the elliptic linear operator acting as

$$
\mathbf{L}_{t} u_{t}(x)=\frac{1}{2} \mathbf{a}_{t}(x) \partial_{x x} u_{t}(x)+\mathbf{b}_{t}(x) \partial_{x} u_{t}(x)+\mathbf{c}_{t}(x) u_{t}(x)
$$

and $\mathbf{G}_{t}$ is the first-order linear operator acting as

$$
\mathbf{G}_{t} u_{t}(x)=\boldsymbol{\sigma}_{t}(x) \partial_{x} u_{t}(x)+\mathbf{g}_{t}(x) u_{t}(x)
$$

The coefficients $(\mathbf{a}, \mathbf{b}, \mathbf{c}, \mathbf{g}, \boldsymbol{\sigma})$ are random fields indexed by $(t, x) \in[0, \infty[\times \mathbb{R}$ and the initial datum $\varphi$ is a random field on $\mathbb{R}$. A classical solution to 4.1 is understood here as a predictable and almost-surely continuous random field $u=u_{t}(x)$ over $\left[0, \infty\left[\times \mathbb{R}\right.\right.$, such that $u_{t} \in C^{2}(\mathbb{R})$ a.s. for any $t>0$ and

$$
u_{t}(x)=\varphi(x)+\int_{0}^{t} \mathbf{L}_{\tau} u_{\tau}(x) d \tau+\int_{0}^{t} \mathbf{G}_{\tau} u_{\tau}(x) d W_{\tau}, \quad t \geq 0, x \in \mathbb{R}
$$

There is a vast literature on stochastic SPDEs and problems of the form 4.1), under suitable measurability, regularity and boundedness assumptions on the coefficients and on the initial datum: see, for instance, [15], [22], 9], [26] and the references therein.

Note that, in analogy with deterministic PDEs, the solution of the Cauchy problem 4.1) can be written, in some cases, as a convolution of the initial datum with a stochastic fundamental solution $p(t, x ; 0, \xi)$, i.e.

$$
\left.u_{t}(x)=\int_{\mathbb{R}} p(t, x ; 0, \xi) \varphi(\xi) d \xi, \quad(t, x) \in\right] 0, \infty[\times \mathbb{R}
$$

with $p(t, x ; 0, \xi)$ being a random field that solves the SPDE in 4.1) with respect to the variables $(t, x)$ and which approximates a Dirac delta centered at $\xi$ as $t$ approaches 0 .

### 4.2 Finite-difference Magnus scheme

We employ the stochastic Magnus expansion to develop an approximation scheme for the Cauchy problem (4.1). Our goal here is only to hint at the possibility that the stochastic Magnus expansion is a useful tool for the numerical solution of SPDEs. Therefore, we keep the exposition at a heuristic level and postpone the rigorous study of the problem for further research.

The idea is to apply finite-difference space-discretization for the operators $\mathbf{L}$ and $\mathbf{G}$, and then Magnus expansion to solve the resulting linear (matrix-valued) Itô SDE. We fix a bounded interval ] $a, b[$ and use the following notation: for a given $d \in \mathbb{N}$, we denote by $\varsigma_{d}$ a mesh of $d+2$ equidistant points in $] a, b[$, i.e.

$$
\varsigma_{d}=\left\{x_{i}^{d} \mid x_{i}^{d}=a+i h, i=0, \ldots, d+1\right\}, \quad h:=\frac{b-a}{d+1}
$$

and for any random field $\mathbf{f}(x), x \in \mathbb{R}$, we denote by $\mathbf{f}^{d}=\left(\mathbf{f}_{0}^{d}, \ldots, \mathbf{f}_{d+1}^{d}\right)$ the random vector whose components correspond to $\mathbf{f}$ evaluated at the points of the mesh, namely

$$
\mathbf{f}_{i}^{d}=\mathbf{f}\left(x_{i}^{d}\right), \quad i=0, \ldots, d+1
$$

Following the classical centered finite-difference discretization, we approximate the spatial derivatives in each point as

$$
\partial_{x} u_{t}\left(x_{i}^{d}\right) \approx \frac{u_{t, i}^{d}-u_{t, i-1}^{d}}{h}, \quad \partial_{x x} u_{t}\left(x_{i}^{d}\right) \approx \frac{u_{t, i+1}^{d}-2 u_{t, i}^{d}+u_{t, i-1}^{d}}{h^{2}}, \quad i=1, \ldots, d
$$

to obtain the system of Itô SDEs

$$
\left\{\begin{array}{l}
d u_{t, i}^{d}=\left(\mathbf{L}_{t}^{d} u_{t}^{d}\right)_{i} d t+\left(\mathbf{G}_{t}^{d} u_{t}^{d}\right)_{i} d W_{t}  \tag{4.2}\\
u_{0, i}^{d}=\varphi_{i}^{d}
\end{array}\right.
$$

for $i=1, \cdots, d$, where $\mathbf{L}_{t}^{d}$ and $\mathbf{G}_{t}^{d}$ are now the operators acting as

$$
\begin{aligned}
\left(\mathbf{L}_{t}^{d} u_{t}^{d}\right)_{i} & =\frac{1}{2} \mathbf{a}_{t, i}^{d} \frac{u_{i+1}^{d}-2 u_{i}^{d}+u_{i-1}^{d}}{h^{2}}+\mathbf{b}_{t, i}^{d} \frac{u_{i}^{d}-u_{i-1}^{d}}{h}+\mathbf{c}_{t, i}^{d} u_{i}^{d} \\
\left(\mathbf{G}_{t}^{d} u_{t}^{d}\right)_{i} & =\boldsymbol{\sigma}_{t, i}^{d} \frac{u_{i}^{d}-u_{i-1}^{d}}{h}+\mathbf{g}_{t, i}^{d} u_{i}^{d}
\end{aligned}
$$

By imposing some boundary conditions, for instance

$$
\begin{equation*}
u_{t, 0}^{d}=u_{t, d+1}^{d}=0, \quad t>0 \tag{4.3}
\end{equation*}
$$

the system of SDEs 4.2 can be cast in the framework of the previous section. More precisely, under condition 4.3 , system 4.2 is equivalent to

$$
\left\{\begin{array}{l}
d \bar{u}_{t}^{d}=A_{t}^{d} \bar{u}_{t}^{d} d t+B_{t}^{d} \bar{u}_{t}^{d} d W_{t}  \tag{4.4}\\
\bar{u}_{0}^{d}=\bar{\varphi}^{d}
\end{array}\right.
$$

where we set

$$
\bar{u}_{t}^{d}=\left(u_{t, 1}^{d}, \ldots, u_{t, d}^{d}\right), \quad \bar{\varphi}^{d}=\left(\varphi_{1}^{d}, \ldots, \varphi_{d}^{d}\right),
$$

and $A_{t}^{d}, B_{t}^{d}$ are the random tridiagonal $(d \times d)$-matrices given by

$$
A_{t}^{d}=\left(\begin{array}{cccc}
-\frac{\mathbf{a}_{t, 1}^{d}}{h^{2}}+\frac{\mathbf{b}_{t, 1}^{d}}{h}+\mathbf{c}_{t, 1}^{d} & \frac{1}{2} \frac{\mathbf{a}_{t, 1}^{d}}{h^{2}} & \cdots & 0 \\
\frac{1}{2} \frac{\mathbf{a}_{t, 2}^{d}}{h^{2}}-\frac{\mathbf{b}_{t, 2}^{d}}{h} & -\frac{\mathbf{a}_{t, 2}^{d}}{h^{2}}+\frac{\mathbf{b}_{t, 2}^{d}}{h}+\mathbf{c}_{t, 2}^{d} & \ddots & \vdots \\
\vdots & \ddots & \ddots & \\
0 & \cdots & \frac{1}{2} \frac{\mathbf{a}_{t, d}^{d}}{h^{2}}-\frac{\mathbf{b}_{t, d}^{d}}{h} & -\frac{\mathbf{a}_{t, d}^{d}}{h^{2}}+\frac{\mathbf{b}_{t, d}^{d}}{h}+\mathbf{c}_{t, d}^{d}
\end{array}\right)
$$

$$
B_{t}^{d}=\left(\begin{array}{cccc}
\frac{\sigma_{t, 1}^{d}}{h}+\mathbf{g}_{t, 1}^{d} & 0 & \cdots & 0 \\
-\frac{\sigma_{t, 2}^{d}}{h} & +\frac{\sigma_{t, 2}^{d}}{h}+\mathbf{g}_{t, 2}^{d} & \ddots & \vdots \\
\vdots & \ddots & \ddots & 0 \\
0 & \cdots & -\frac{\sigma_{t, d}^{d}}{h} & +\frac{\sigma_{t, d}^{d}}{h}+\mathbf{g}_{t, d}^{d}
\end{array}\right)
$$

Now, the solution to (4.4) can be written as

$$
\bar{u}_{t}^{d}=X_{t}^{d} \bar{u}_{0}^{d}, \quad t \geq 0,
$$

where $X^{d}$ is in turn the solution to the $\mathcal{M}^{d \times d}$-valued Itô SDE

$$
\left\{\begin{array}{l}
d X_{t}^{d}=A_{t}^{d} X_{t}^{d} d t+B_{t}^{d} X_{t}^{d} d W_{t},  \tag{4.5}\\
X_{0}^{d}=I_{d} .
\end{array}\right.
$$

Remark 4.1. The components of $X^{d}$ can be regarded as approximations of the integrals of the fundamental solution of the SPDE in 4.1], when it exists, on each sub-interval $\left[\frac{1}{2}\left(x_{j-1}^{d}+x_{j}^{d}\right), \frac{1}{2}\left(x_{j}^{d}+x_{j+1}^{d}\right)\right]$, namely

$$
\begin{equation*}
\left(X_{t}^{d}\right)_{i, j} \approx \int_{\frac{1}{2}\left(x_{j-1}^{d}+x_{j}^{d}\right)}^{\frac{1}{2}\left(x_{j}^{d}+x_{j+1}^{d}\right)} p\left(t, x_{i}^{d} ; 0, \xi\right) d \xi=:\left(\mathcal{I}_{t}^{d}\right)_{i, j}, \quad i, j=1, \ldots, d . \tag{4.6}
\end{equation*}
$$

### 4.3 A numerical test

In this section we consider a special case of (4.1) with $\mathbf{a}_{t} \equiv \mathbf{a}>0, \mathbf{b}, \mathbf{c}, \mathbf{g} \equiv 0$ and $\sigma_{t} \equiv \sigma>0$. Hence, we consider the stochastic heat equation

$$
d u_{t}=\frac{\mathbf{a}}{2} \partial_{x x} u_{t}(x) d t+\boldsymbol{\sigma} \partial_{x} u_{t}(x) d W_{t}, \quad t>0, x \in \mathbb{R}
$$

with $\mathbf{a}>\boldsymbol{\sigma}^{2}$, whose stochastic fundamental solution is given explicitly by

$$
\begin{equation*}
p(t, x ; 0, \xi):=\frac{1}{\sqrt{2 \pi\left(\mathbf{a}-\boldsymbol{\sigma}^{2}\right) t}} \exp \left(-\frac{\left(x+\boldsymbol{\sigma} W_{t}-\xi\right)^{2}}{2\left(\mathbf{a}-\boldsymbol{\sigma}^{2}\right) t}\right), \quad t>0, x, \xi \in \mathbb{R} . \tag{4.7}
\end{equation*}
$$

The matrices $A_{t}^{d}$ and $B_{t}^{d}$ in (4.5) now read as

$$
A_{t}^{d} \equiv \frac{\mathbf{a}}{h^{2}}\left(\begin{array}{cccc}
-1 & \frac{1}{2} & \cdots & 0 \\
\frac{1}{2} & -1 & \ddots & \vdots \\
\vdots & \ddots & \ddots & \frac{1}{2} \\
0 & \cdots & \frac{1}{2} & -1
\end{array}\right), \quad B_{t}^{d} \equiv \frac{\sigma}{h}\left(\begin{array}{cccc}
1 & 0 & \cdots & 0 \\
-1 & 1 & \ddots & \vdots \\
\vdots & \ddots & \ddots & 0 \\
0 & \cdots & -1 & 1
\end{array}\right),
$$

and do not commute.
In the next numerical test we compare the approximate solutions to 4.5), obtained with the stochastic Magnus expansion, with the $\mathcal{M}^{d \times d}$-valued stochastic process $\mathcal{I}^{d}$, whose components are given by the integral

| Method | $t=0.1$ | $t=0.2$ | $t=0.3$ | $t=0.4$ | $t=0.5$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| m 1 | $7.535 \%$ | $5.497 \%$ | $4.663 \%$ | $4.790 \%$ | $4.539 \%$ |
| m 3 | $8.825 \%$ | $6.716 \%$ | $5.445 \%$ | $5.142 \%$ | $4.329 \%$ |

Table 4: SPDE. Values of $\mathbb{E}\left[\operatorname{Err}_{t}^{d}\right]$ (in percentage) for m 1 and m 3 , with $d=50$, obtained with 50 independent samples. Parameters as in 4.9.

| Method | $t=0.1$ | $t=0.2$ | $t=0.3$ | $t=0.4$ | $t=0.5$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| m 1 | $4.793 \%$ | $4.625 \%$ | $4.703 \%$ | $4.824 \%$ | $4.665 \%$ |
| m 3 | $4.464 \%$ | $3.216 \%$ | $2.482 \%$ | $2.151 \%$ | $1.949 \%$ |

Table 5: SPDE. Values of $\mathbb{E}\left[\operatorname{Err}_{t}^{d}\right]$ (in percentage) for m 1 and m 3 , with $d=100$, obtained with 50 independent samples. Parameters as in 4.9.

| Method | $t=0.1$ | $t=0.2$ | $t=0.3$ | $t=0.4$ | $t=0.5$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| m 1 | $4.980 \%$ | $5.152 \%$ | $5.056 \%$ | $5.064 \%$ | $5.073 \%$ |
| m 3 | $2.352 \%$ | $1.813 \%$ | $1.304 \%$ | $1.045 \%$ | $0.941 \%$ |

Table 6: SPDE. Values of $\mathbb{E}\left[\operatorname{Err}_{t}^{d}\right]$ (in percentage) for m 1 and m 3 , with $d=200$, obtained with 50 independent samples. Parameters as in 4.9.
in 4.6) with $p$ as in 4.7. In doing this, we shall keep in mind that the difference between the latter quantities can be decomposed into two errors, namely: the one between $X^{d}$ and its approximation, and the one between $X^{d}$ and $\mathcal{I}^{d}$. In turn, the latter is the result of both space-discretization and the error that stems by imposing null boundary conditions (see 4.3). In particular, this last error cannot be reduced by refining the spacegrid. Therefore, the analysis should be restricted to the "central" components of $\mathcal{I}^{d}$, namely those which do not depend on the values of the fundamental solution in the vicinity of the boundary $\{a, b\}$. This motivates the definition that follows. For a given $\kappa \in \mathbb{N}$ with $\kappa<d$, and a given approximation $X^{d, \text { app }}$ of $X^{d}$, we define the process

$$
\begin{equation*}
\operatorname{Err}_{t}^{d}:=\frac{\left|\widetilde{\mathcal{I}}_{t}^{d, \kappa}-\widetilde{X}_{t}^{d, \kappa, \mathrm{app}}\right|}{\left|\widetilde{\mathcal{I}}_{t}^{d, \kappa}\right|}, \quad t>0 \tag{4.8}
\end{equation*}
$$

where $\widetilde{\mathcal{I}}_{t}^{d, \kappa}$ and $\widetilde{X}_{t}^{d, \kappa, \text { app }}$ are the projections on $\mathcal{M}^{\kappa \times d}$ obtained by selecting the central $\kappa$ rows of $\mathcal{I}_{t}^{d}$ and $X_{t}^{d, \text { app }}$, respectively. The matrix norm above is the Frobenius norm. The role of $X^{d, \text { app }}$ will be played by the time-discretization of the truncated Magnus expansion 1.8 - 1.9 . In particular, we will denote by m1 and $m 3$ the discretized first and third-order Magnus expansions of $X^{d}$, respectively. We will not consider the second-order Magnus approximation m 2 as it appears less stable than the others. Note that, being $A_{t}^{d}$ and $B_{t}^{d}$ constant matrices, the first three terms of the Magnus expansion are given explicitly by 2.21 .

In the numerical experiments we set

$$
\begin{equation*}
a=-2, b=2, \quad \mathbf{a}=0.2, \quad \boldsymbol{\sigma}=0.15 \tag{4.9}
\end{equation*}
$$

Setting the parameter $\kappa$ in 4.8, which determines the number of rows that are taken into account to asses the error, as $\kappa=\lfloor d / 2\rfloor$, we study the expectation of $\operatorname{Err}_{t}^{d}$ up to $t=0.5$. Such choice for $\kappa$ and $t$ allows us to study the error in a region that is suitably away from the boundary. Indeed, choosing $\kappa$ as above implies
$x_{i}^{d}$ in (4.6) ranging roughly from -1 to 1 . On the other hand, the standard-deviation parameter associated to the Gaussian density 4.7) at $t=0.5$ is roughly 0.30 , while the mean parameter is $0.15 \times W_{0.5}$, whose standard deviation is in turn roughly 0.10 . Therefore, both $\left(\widetilde{\mathcal{I}}_{t}^{d, \kappa}\right)_{i, 1}$ and $\left(\widetilde{\mathcal{I}}_{t}^{d, \kappa}\right)_{i, d}$ are likely to be very close to zero, thus meeting the null boundary condition implied by 4.3.

In Tables 4. 5, 6, we report the approximate values of $\mathbb{E}\left[\operatorname{Err}_{t}^{d}\right]$ for $d=50,100$ and 200, respectively. These were obtained via simulation of 50 trajectories of $W$ with $10^{4}$ time steps.

## 5 Conclusion

In this paper we derive the stochastic Magnus expansion for Itô SDEs and prove its convergence up to a positive stopping time $\tau$. Our main result, Theorem 1.1, also provides an asymptotic estimate for $\tau$.

Additionally, we test the Magnus expansion for the numerical solution of stochastic differential equations, confirming its potential as a highly parallelizable and accurate method. The Magnus expansion also provides a novel method for solving SPDEs for which we present some preliminary tests. A complete theoretical error analysis is left to future research, as well as possible extensions to more general classes of stochastic processes including jumps and mean-field interactions.

From a numerical perspective, some of the main advantages and drawbacks of the Magnus expansion are summarized in the following table:

| Pros | Cons |
| :--- | :--- |
| Calculation of the logarithm is very fast | Calculation of the exponential can be slow with <br> standard libraries |
| Can be used to evaluate multiple times very fast | As its deterministic counterpart, it works only up to <br> a finite time |
| Highly parallelizable and GPU applicable | Memory demanding in high dimension |

## A Derivatives of matrix exponentials

In the next lemma we provide explicit expressions for the first and second order differentials of the exponential $\operatorname{map} \mathcal{M}^{d \times d} \ni M \mapsto e^{M}$. We recall that this map is smooth and in particular, it is continuously twice differentiable.

Lemma A.1. For any $\Sigma \in \mathcal{M}^{d \times d}$, the first and the second order differentials at $\Sigma$ of the exponential map $\mathcal{M}^{d \times d} \ni M \mapsto e^{M}$ are given by

$$
\begin{align*}
M & \mapsto \mathcal{L}_{\Sigma}(M) e^{\Sigma}=e^{\Sigma} \mathcal{L}_{-\Sigma}(M), & & M \in \mathcal{M}^{d \times d}  \tag{A.1}\\
(M, N) & \mapsto \mathcal{Q}_{\Sigma}(M, N) e^{\Sigma}, & & M, N \in \mathcal{M}^{d \times d}
\end{align*}
$$

where $\mathcal{L}_{\Sigma}$ and $\mathcal{Q}_{\Sigma}$ are the linear and (symmetric) bi-linear operators as defined in 2.2 - 2.3 .
Proof. The first part of the statement, concerning the first order differential, is a classical result; its proof can be found in [3, Lemma 2] among other references.

We prove the second part. Fix $M \in \mathcal{M}^{d \times d}$ and denote by $\partial_{M} e^{\Sigma}$ the first order directional derivative of $e^{\Sigma}$ w.r.t. $M$, i.e.

$$
\partial_{M} e^{\Sigma}:=\left.\frac{d}{d t} e^{\Sigma+t M}\right|_{t=0}
$$

By the first part, we have

$$
\partial_{M} e^{\Sigma}=\mathcal{L}_{\Sigma}(M) e^{\Sigma}, \quad \Sigma \in \mathcal{M}^{d \times d}
$$

We now show that, for any $M, N \in \mathcal{M}^{d \times d}$, the second order directional derivative

$$
\partial_{N, M} e^{\Sigma}:=\left.\frac{d}{d t} \partial_{M} e^{\Sigma+t N}\right|_{t=0}
$$

is given by

$$
\begin{equation*}
\partial_{N, M} e^{\Sigma}=Q_{\Sigma}(N, M) e^{\Sigma}, \quad \Sigma \in \mathcal{M}^{d \times d} \tag{A.2}
\end{equation*}
$$

We have

$$
\begin{align*}
\frac{d}{d t} \partial_{M} e^{\Sigma+t N} & =\frac{d}{d t}\left(\mathcal{L}_{\Sigma+t N}(M) e^{\Sigma+t N}\right) \\
& =\mathcal{L}_{\Sigma+t N}(M) \mathcal{L}_{\Sigma+t N}(N) e^{\Sigma+t N}+\left(\frac{d}{d t} \mathcal{L}_{\Sigma+t N}(M)\right) e^{\Sigma+t N} \tag{A.3}
\end{align*}
$$

We use the definition 2.2 and exchange the differentiation and integration signs to obtain

$$
\frac{d}{d t} \mathcal{L}_{\Sigma+t N}(M)=\int_{0}^{1} \frac{d}{d t} e^{\operatorname{ad}_{\tau(\Sigma+t N)}}(M) d \tau
$$

(by 2.1) )

$$
\begin{aligned}
& =\int_{0}^{1} \frac{d}{d t}\left(e^{\tau(\Sigma+t N)} M e^{-\tau(\Sigma+t N)}\right) d \tau \\
& =\int_{0}^{1}\left(\frac{d}{d t} e^{\tau(\Sigma+t N)}\right) M e^{-\tau(\Sigma+t N)} d \tau+\int_{0}^{1} e^{\tau(\Sigma+t N)} M \frac{d}{d t} e^{-\tau(\Sigma+t N)} d \tau
\end{aligned}
$$

(by employing the two expressions in A.1 for the first-order differential)

$$
\begin{aligned}
= & \int_{0}^{1} \tau \mathcal{L}_{\tau(\Sigma+t N)}(N) e^{\tau(\Sigma+t N)} M e^{-\tau(\Sigma+t N)} d \tau \\
& -\int_{0}^{1} \tau e^{\tau(\Sigma+t N)} M e^{-\tau(\Sigma+t N)} \mathcal{L}_{\tau(\Sigma+t N)}(N) d \tau \\
= & \int_{0}^{1} \tau\left[\mathcal{L}_{\tau(\Sigma+t N)}(N), e^{\operatorname{ad}_{\tau(\Sigma+t N)}}(M)\right] d \tau
\end{aligned}
$$

This, together with A.3), proves A.2.
To conclude, we prove equality (2.4). It is enough to observe that

$$
\begin{aligned}
\int_{0}^{1} \tau \mathcal{L}_{\tau \Sigma}(N) e^{\operatorname{ad}_{\tau \Sigma}(M) d \tau} & =\int_{0}^{1} \tau \sum_{n=0}^{\infty} \sum_{m=0}^{\infty} \frac{\operatorname{ad}_{\tau \Sigma}^{n}(N)}{(n+1)!} \frac{\operatorname{ad}_{\tau \Sigma}^{m}(M)}{m!} d \tau \\
& =\int_{0}^{1} \sum_{n=0}^{\infty} \sum_{m=0}^{\infty} \frac{\operatorname{ad}_{\Sigma}^{n}(N)}{(n+1)!} \frac{\operatorname{ad}_{\Sigma}^{m}(M)}{m!} \tau^{n+m+1} d \tau \\
& =\sum_{n=0}^{\infty} \sum_{m=0}^{\infty} \frac{\operatorname{ad}_{\Sigma}^{n}(N) \operatorname{ad}_{\Sigma}^{m}(M)}{(n+m+2)(n+1)!m!}
\end{aligned}
$$

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