An efficient algorithm for computing the Baker–Campbell–Hausdorff series and some of its applications

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We provide a new algorithm for generating the Baker–Campbell–Hausdorff (BCH) series $Z = \log(e^X e^Y)$ in an arbitrary generalized Hall basis of the free Lie algebra $\mathcal{L}(X, Y)$ generated by X and Y. It is based on the close relationship of $\mathcal{L}(X, Y)$ with a Lie algebraic structure of labeled rooted trees. With this algorithm, the computation of the BCH series up to degree of 20 [111 013 independent elements in $\mathcal{L}(X, Y)$] takes less than 15 min on a personal computer and requires 1.5 Gbytes of memory. We also address the issue of the convergence of the series, providing an optimal convergence domain when X and Y are real or complex matrices. © *2009 American Institute of Physics.* [DOI: 10.1063/1.3078418]

I. INTRODUCTION

The Baker–Campbell–Hausdorff (BCH) formula deals with the expansion of Z in $e^X e^Y = e^Z$ in terms of nested commutators of X and Y when they are assumed to be noncommuting operators. If we introduce the formal series for the exponential function

$$e^X e^Y = \sum_{p,q=0}^{\infty} \frac{1}{p! q!} X^p Y^q \tag{1.1}$$

and substitute this series in the formal series defining the logarithm function

$$\log Z = \sum_{k=1}^{\infty} \frac{(-1)^{k-1}}{k} (Z - 1)^k,$$

one obtains

$$\log(e^X e^Y) = \sum_{k=1}^{\infty} \frac{(-1)^{k-1}}{k} \sum \frac{X^{p_1} Y^{q_1} \cdots X^{p_k} Y^{q_k}}{p_1! q_1! \cdots p_k! q_k!},$$

where the inner summation extends over all non-negative integers $p_1, q_1, \ldots, p_k, q_k$ for which $p_i + q_i > 0$ ($i = 1, 2, \ldots, k$). Gathering together the terms for which $p_1 + q_1 + p_2 + q_2 + \cdots + p_k + q_k = m$, we can write

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$$Z = \log(e^X e^Y) = \sum_{m=1}^{\infty} P_m(X,Y), \tag{1.2}$$

where $P_m(X,Y)$ is a homogeneous polynomial of degree m in the noncommuting variables X and Y. Campbell,[8] Baker,[2] and Hausdorff[17] addressed the question whether Z can be represented as a series of nested commutators of X and Y, without producing a general formula. We recall here that the commutator $[X,Y]$ is defined as $XY - YX$. It was Dynkin[12] who finally derived an explicit formula for Z as

$$Z = \sum_{k=1}^{\infty} \sum_{p_i, q_i} \frac{(-1)^{k-1}}{k} \frac{[X^{p_1} Y^{q_1} \cdots X^{p_k} Y^{q_k}]}{(\sum_{i=1}^{k}(p_i + q_i))p_1! q_1! \cdots p_k! q_k!}. \tag{1.3}$$

Here the inner summation is taken over all non-negative integers $p_1, q_1, \ldots, p_k, q_k$ such that $p_1 + q_1 > 0, \ldots, p_k + q_k > 0$ and $[X^{p_1} Y^{q_1} \cdots X^{p_k} Y^{q_k}]$ denotes the right nested commutator based on the *word* $X^{p_1} Y^{q_1} \cdots X^{p_k} Y^{q_k}$. Expression (1.3) is known, for obvious reasons, as the BCH *series in the Dynkin form*. By rearranging terms, it is clear that Z can be written as

$$Z = \log(e^X e^Y) = X + Y + \sum_{m=2}^{\infty} Z_m, \tag{1.4}$$

with $Z_m(X,Y)$ a homogeneous Lie polynomial in X and Y of degree m, i.e., it is a $\mathbb{Q}$-linear combination of commutators of the form $[V_1, [V_2, \ldots, [V_{m-1}, V_m] \cdots]]$ with $V_i \in \{X, Y\}$ for $1 \leq i \leq m$. The first terms read explicitly

$$Z_2 = \tfrac{1}{2}[X,Y],$$

$$Z_3 = \tfrac{1}{12}[X,[X,Y]] - \tfrac{1}{12}[Y,[X,Y]],$$

$$Z_4 = \tfrac{1}{24}[X,[Y,[Y,X]]].$$

The expression $e^X e^Y = e^Z$ is then called the BCH *formula*, although other different labels (e.g., Campbell–Baker–Hausdorff, Baker–Hausdorff, Campbell–Hausdorff) are commonly attached to it in the literature. The formula (1.3) is certainly awkward to use due to the complexity of the sums involved. Notice, in particular, that different choices of p_i, q_i, k in (1.3) may lead to terms in the same commutator. Thus, for instance, $[X^3 Y^1] = [X^1 Y^0 X^2 Y^1] = [X,[X,[X,Y]]]$. An additional difficulty arises from the fact that not all the commutators are independent due to the Jacobi identity,[47]

$$[X_1, [X_2, X_3]] + [X_2, [X_3, X_1]] + [X_3, [X_1, X_2]] = 0.$$

The BCH formula plays a fundamental role in many fields of mathematics (theory of linear differential equations,[26] Lie groups,[14] numerical analysis[16]), theoretical physics (perturbation theory,[10] quantum mechanics,[49] statistical mechanics,[24,50] quantum computing[40]), and control theory (analysis and design of nonlinear control laws, nonlinear filters, stabilization of rigid bodies[46]). In particular, in the theory of Lie groups, with this formula one can explicitly write the operation of multiplication in a Lie group in canonical coordinates in terms of the Lie bracket operation in its tangent algebra and also prove the existence of a local Lie group with a given Lie algebra.[14]

Also in the numerical treatment of differential equations on manifolds,[19,16] the BCH formula is quite useful. If $\mathcal{M}$ is a smooth manifold and $\mathcal{X}(\mathcal{M})$ denotes the linear space of smooth vector fields on $\mathcal{M}$, then a Lie algebra structure is established in $\mathcal{X}(\mathcal{M})$ by using the Lie bracket $[X,Y]$ of fields X and $Y \in \mathcal{X}(\mathcal{M})$.[47] The flow of a vector field $X \in \mathcal{X}(\mathcal{M})$ is a mapping $\exp(X)$ defined through the solution of the differential equation

$$\frac{du}{dt} = X(u), \quad u(0) = q \in \mathcal{M} \tag{1.5}$$

as $\exp(tX)(q) = u(t)$. Many numerical methods used to approximately solve Eq. (1.5) are based on compositions of maps that are flows of vector fields.[16] To be more specific, suppose the vector field X can be split as $X = A + B$ and that the flows corresponding to $A(u)$ and $B(u)$ can be explicitly obtained. Then one may consider an approximation of the form $\Psi_h \equiv \exp(ha_1 A)\exp(hb_1 B)\cdots\exp(ha_k A)\exp(hb_k B)$ for the exact flow $\exp(h(A+B))$ of (1.5) after a time step h. The idea now is to obtain the conditions to be satisfied by the coefficients a_i, b_i so that $\Psi_h(q) = u(h) + \mathcal{O}(h^{p+1})$ as $h \to 0$, and this can be done by applying the BCH formula in sequence to the expression of Ψ up to the degree required by the order of approximation p.[27] This task can be carried out quite easily provided one has explicit expressions of Z_m implemented in a symbolic algebra package.[23,46]

In addition to the Dynkin form (1.3), there are other standard procedures to construct explicitly the BCH series. Recall that the free Lie algebra $\mathcal{L}(X, Y)$ generated by the symbols X and Y can be considered as a subspace (the subspace of Lie polynomials) of the vector space spanned by the words w in the symbols X and Y, i.e., $w = a_1 a_2 \cdots a_m$, each a_i being X or Y. Thus, the BCH series admits the explicit associative presentation

$$Z = X + Y + \sum_{m=2}^{\infty} \sum_{w, |w|=m} g_w w, \tag{1.6}$$

in which g_w is a rational coefficient and the inner sum is taken over all words w with length $|w| = m$. Here the length of w is the number of letters it contains. The coefficients can be computed with a procedure based on a family of recursively computable polynomials.[13]

Although the terms in Eq. (1.6) are expressed as linear combinations of individual words (which are not Lie polynomials), by virtue of the Dynkin–Specht–Wever theorem,[21] Z can be written as

$$Z = X + Y + \sum_{m=2}^{\infty} \frac{1}{m} \sum_{w, |w|=m} g_w [w], \tag{1.7}$$

that is, the individual terms are the same as in the associative series (1.6) except that the word $w = a_1 a_2 \cdots a_m$ is replaced with the right nested commutator $[w] = [a_1, [a_2, \ldots, [a_{m-1}, a_m] \cdots]]$ and the coefficient g_w is divided by the word length m.[42] This gives explicit expressions of the terms Z_m in the BCH series (1.4) as a linear combination of nested commutators of homogeneous degree, that is, as a linear combination of elements of the homogeneous subspace $\mathcal{L}(X, Y)_m$ of degre m of the free Lie algebra $\mathcal{L}(X, Y)$. However, it should be stressed that the set of nested commutators $[w]$ for words w of length m is not a basis of the homogeneous subspace $\mathcal{L}(X, Y)_m$.

By introducing a parameter τ and differentiating with respect to τ the power series $\sum_{m \geq 1} \tau^m Z_m = \log(\exp(\tau X)\exp(\tau Y))$, the following recursion formula is derived in Ref. 47:

$$Z_1 = X + Y,$$

$$mZ_m = \frac{1}{2}[X - Y, Z_{m-1}] + \sum_{p=1}^{[(m-1)/2]} \frac{B_{2p}}{(2p)!}(\text{ad}_Z^{2p}(X+Y))_m, \quad m \geq 1. \tag{1.8}$$

Here $Z = \sum_{m \geq 1} Z_m$, $\text{ad}_Z^k(X+Y) = [Z, \text{ad}_Z^{k-1}(X+Y)]$, the B_j stand for the Bernoulli numbers,[1] and $(\text{ad}_Z^{2p}(X+Y))_m$ denotes the projection of $\text{ad}_Z^{2p}(X+Y)$ onto the homogeneous subspace $\mathcal{L}(X, Y)_m$, which can be written in terms of $Z_1, Z_2, Z_3, \ldots$ as

$$(\mathrm{ad}_Z^{2p}(X+Y))_m = \sum_{\substack{k_1+\cdots+k_{2p}=m-1 \\ k_1 \geqslant 1,\ldots,k_{2p} \geqslant 1}} [Z_{k_1},[\cdots[Z_{k_{2p}},X+Y]\cdots]].$$

Explicit formulas (1.3) and (1.7), as well as recursion (1.8) can be used in principle to construct the BCH series up to arbitrary degree in terms of commutators. As a matter of fact, several systematic computations of the series have been carried out along the years, starting with the work of Richtmyer and Greenspan in 1965,[37] where results up to degree of 8 are reported. Later on, Newman and Thompson obtained the coefficients g_w in (1.7) up to words of length of 20,[32] Bose[6] constructed an algorithm to compute directly the coefficient of a given commutator in the Dynkin presentation (1.3) and Oteo[33] and Kolsrud[22] presented a simplified expression of (1.3) in terms of right nested commutators up to degrees of 8 and 9, respectively. More recently, Reinsch[35] proposed a matrix operation procedure for calculating the polynomials $P_m(X,Y)$ in (1.2) which can be easily implemented in any symbolic algebra package. Again, the Dynkin–Specht–Wever has to be used to write the resulting expressions in terms of commutators.

As mentioned before, all of these procedures exhibit a key limitation, however: the iterated commutators are not all linearly independent due to the Jacobi identity (and other identities involving nested commutators of higher degree which are originated by it[33]). In other words, they do not provide expressions directly in terms of a basis of the free Lie algebra $\mathcal{L}(X,Y)$. This is required, for instance, in applications of the BCH formula in the numerical integration of ordinary differential equations or when one wants to study specific features of the series, such as the distribution of the coefficients and other combinatorial properties.[32]

Of course, it is always possible to express the resulting formulas in terms of a basis of $\mathcal{L}(X,Y)$ but this rewriting process is very time consuming and requires a good deal of memory resources. In practice, going beyond degree $m=11$ constitutes a difficult task indeed,[28,23,46] since the number of terms involved in the series grows, in general, as the dimension c_m of the homogeneous subspace $\mathcal{L}(X,Y)_m$. As is well known, c_m is given by the Witt formula,[7] so that $c_m = \mathcal{O}(2^m/m)$.

Our goal is then to express the BCH series as

$$Z = \log(\exp(X)\exp(Y)) = \sum_{i \geqslant 1} z_i E_i, \tag{1.9}$$

where $z_i \in \mathbb{Q}$ ($i \geqslant 1$) and $\{E_i : i = 1,2,3,\ldots\}$ is a basis of $\mathcal{L}(X,Y)$ whose elements are of the form

$$E_1 = X, \quad E_2 = Y, \quad \text{and} \quad E_i = [E_{i'},E_{i''}], \quad i \geqslant 3, \tag{1.10}$$

for appropriate values of the integers $i',i'' < i$ (for $i=3,4,\ldots$). Clearly, each E_i in (1.10) is a homogeneous Lie polynomial of degree $|i|$, where

$$|1| = |2| = 1 \quad \text{and} \quad |i| = |i'| + |i''| \quad \text{for } i \geqslant 3. \tag{1.11}$$

We will focus on a general class of bases of the free Lie algebra $\mathcal{L}(X,Y)$, referred to in the current literature as generalized Hall bases and also as Hall–Viennot bases.[36,48] These include the Lyndon basis[25,48] and different variants of the classical Hall basis (see Ref. 36, for references). Specifically, in this paper we present a new procedure to write the BCH series (1.9) for an arbitrary Hall–Viennot basis. Such an algorithm is based on results obtained in Ref. 30, in particular, those relating a certain Lie algebra structure $\mathfrak{g}$ on rooted trees with the description of a free Lie algebra in terms of a Hall basis. This Lie algebra $\mathfrak{g}$ on rooted trees was first considered in Ref. 11, whereas a closely related Lie algebra on labeled rooted trees was treated in Ref. 15 (see Ref. 18 for the relation of these two Lie algebras and for further references about related algebraic structures on rooted trees).

We have implemented the algorithm in MATHEMATICA (it can also be programmed in FORTRAN or C for more efficiency). The resulting procedure gives the BCH series up to a prescribed degree directly in terms of a Hall–Viennot basis of $\mathcal{L}(X,Y)$. As an illustration, obtaining the series (in the classical basis of P. Hall) up to degree $m=20$ with a personal computer (2.4 GHz Intel Core 2 Duo

processor with 2 Gbytes of random access memory) requires less than 15 min of CPU time and 1.5 Gbytes of memory. The resulting expression has 109 697 nonvanishing coefficients out of 111 013 elements E_i of degree $|i| \leq 20$ in the Hall basis. As far as we know, there are no results up to such a high degree reported in the literature. For comparison with other procedures, the authors of Ref. 46 reported 25 h of CPU time and 17.5 Mbytes with a Pentium III personal computer to achieve degree of 10. By contrast, our algorithm is able to achieve $m = 10$ in 0.058 s and only needs 5.4 Mbytes of computer memory.

In Table III in the Appendix, we give the values of i' and i'' for the elements E_i of degree $|i| \leq 9$ in the Hall basis and their coefficients z_i in the BCH formula (1.9). The elements of the basis are ordered in such a way that $i < j$ if $|i| < |j|$, and the horizontal lines in the table separate elements of different homogeneous degree. Extension of Table III up to terms of degree of 20 is available at the website www.gicas.uji.es/research/bch.html for both the basis of P. Hall and the Lyndon basis. As an example, the last element of degree of 20 in the Hall basis is

$$E_{111013} = [[[[[Y,X],Y],[Y,X]],[[[Y,X],X],[Y,X]]],[[[[Y,X],Y],[Y,X]],[[[[Y,X],Y],Y],Y]]],$$

and the corresponding coefficient in (1.9) reads

$$z_{111\,013} = -\frac{19\,234\,697}{140\,792\,940\,288}.$$

Another central issue addressed in this paper concerns the convergence properties of the BCH series. Suppose we introduce a submultiplicative norm $\|\cdot\|$ such that

$$\|[X,Y]\| \leq \mu \|X\| \|Y\| \tag{1.12}$$

for some $\mu > 0$. Then it is not difficult to show that the series (1.3) is absolutely convergent as long as $\|X\| + \|Y\| < (\log 2)/\mu$.[7,41] As a matter of fact, several improved bounds have been obtained for the different presentations. Thus, in particular, the Lie presentation (1.7) converges absolutely if $\|X\| \leq 1/\mu$ and $\|Y\| \leq 1/\mu$ in a normed Lie algebra $\mathfrak{g}$ with a norm satisfying (1.12),[31,45] whereas in Ref. 3 it has been shown that the series $Z = \Sigma_{m \geq 1} Z_m$ is absolutely convergent for all X, Y such that

$$\mu \|X\| \leq \int_{\mu\|Y\|}^{2\pi} \frac{1}{2 + \frac{t}{2}\left(1 - \cot\left(\frac{1}{2}t\right)\right)} dt \tag{1.13}$$

and the corresponding expression obtained by interchanging in (1.13) X by Y. Moreover, the series diverges, in general, if $\|X\| + \|Y\| \geq \pi$ when $\mu = 2$.[28] Here we provide a generalization of this feature based on the well known Magnus expansion for linear differential equations[26] and also we give a more precise characterization of the convergence domain of the series when X and Y are (real or complex) matrices.

II. AN ALGORITHM FOR COMPUTING THE BCH SERIES BASED ON ROOTED TREES

A. Summary of the procedure

Our starting point is the vector space $\mathfrak{g}$ of maps $\alpha : \mathcal{T} \to \mathbb{R}$, where $\mathcal{T}$ denotes the set of rooted trees with black and white vertices,

$$\mathcal{T} = \left\{ \bullet, \circ, \substack{\bullet\\\bullet}, \ldots, \ldots \right\}.$$

In the combinatorial literature, $\mathcal{T}$ is typically referred to as the set of labeled rooted trees with two labels, "black" and "white." Hereafter, we refer to the elements of $\mathcal{T}$ as bicoloured rooted trees.

The vector space $\mathfrak{g}$ is endowed with a Lie algebra structure by defining the Lie bracket $[\alpha, \beta] \in \mathfrak{g}$ of two arbitrary maps $\alpha, \beta \in \mathfrak{g}$ as follows. For each $u \in \mathcal{T}$,

$$[\alpha,\beta](u) = \sum_{j=1}^{|u|=1} (\alpha(u_{(j)})\beta(u^{(j)}) - \alpha(u^{(j)})\beta(u_{(j)})), \tag{2.1}$$

where $|u|$ denotes the number vertices of u, and each of the pairs of trees $(u_{(j)}, u^{(j)}) \in \mathcal{T} \times \mathcal{T}$, $j = 1, \dots, |u| - 1$, is obtained from u by removing one of the $|u| - 1$ edges of the rooted tree u, the root of $u_{(j)}$ being the original root of u. For instance,

$$[\alpha,\beta](\text{⚇}) = \alpha(\circ)\beta(\bullet) - \alpha(\bullet)\beta(\circ), \qquad [\alpha,\beta](\text{⚇}) = 0,$$

$$[\alpha,\beta](\text{⚇}) = 2(\alpha(\text{⚇})\beta(\bullet) - \alpha(\bullet)\beta(\text{⚇})),$$

$$[\alpha,\beta](\text{⚇}) = \alpha(\text{⚇})\beta(\bullet) + \alpha(\text{⚇})\beta(\circ) - \alpha(\bullet)\beta(\text{⚇}) - \alpha(\circ)\beta(\text{⚇}). \tag{2.2}$$

An important feature of the Lie algebra $\mathfrak{g}$ is that the Lie subalgebra of $\mathfrak{g}$ generated by the maps X, $Y \in \mathfrak{g}$ defined as

$$X(u) = \begin{cases} 1 & \text{if } u = \bullet \\ 0 & \text{if } u \in \mathcal{T} \backslash \{\bullet\} \end{cases}, \qquad Y(u) = \begin{cases} 1 & \text{if } u = \circ \\ 0 & \text{if } u \in \mathcal{T} \backslash \{\circ\} \end{cases}. \tag{2.3}$$

is a free Lie algebra over the set $\{X, Y\}$.[30] In what follows, we denote as $\mathcal{L}(X, Y)$ the Lie subalgebra of $\mathfrak{g}$ generated by the maps X and Y.

It has also been shown in Ref. 30 that for each particular Hall–Viennot basis $\{E_i : i = 1, 2, 3, \dots\}$, [whose elements are given by (1.10) for appropriate values of $i', i'' < i, i = 3, 4, \dots$, and X and Y given by (2.3)] *one can associate a bicoloured rooted tree u_i with each element E_i such that, for any map $\alpha \in \mathcal{L}(X, Y)$,*

$$\alpha = \sum_{i \geq 1} \frac{\alpha(u_i)}{\sigma(u_i)} E_i, \tag{2.4}$$

where for each i, $\sigma(u_i)$ is certain positive integer associated with the bicolored rooted tree u_i (the number of symmetries of u_i, that we call *symmetry number of u_i*). For instance, the bicolored rooted trees u_i and their symmetry numbers $\sigma(u_i)$ associated with the elements E_i (of degree $|i| \leq 5$) of the Hall basis used in this work are displayed in Table I.

As in Sec. I, we denote by $\mathcal{L}(X, Y)_n$ ($n \geq 1$) the homogeneous subspace of $\mathcal{L}(X, Y)$ of degree n (whence admiting $\{E_i : |i| = n\}$ as a basis). It can be seen[30] that if $\alpha \in \mathcal{L}(X, Y)$, then its projection α_n to the homogeneous subspace $\mathcal{L}(X, Y)_n$ is given by

$$\alpha_n(u) = \begin{cases} \alpha(u) & \text{if } |u| = n \\ 0 & \text{otherwise} \end{cases} \tag{2.5}$$

for each $u \in \mathcal{T}$.

We also use the notation $\overline{\mathcal{L}(X, Y)}$ for the Lie algebra of Lie series, that is, series of the form

$$\alpha = \alpha_1 + \alpha_2 + \alpha_3 + \cdots, \quad \text{where } \alpha_n \in \mathcal{L}(X, Y)_n.$$

Notice that in this setting, a Lie series $\alpha \in \overline{\mathcal{L}(X, Y)}$ is a map $\alpha : \mathcal{T} \to \mathbb{R}$ satisfying that, for each $n \geq 1$, the map α_n given by (2.5) belongs to $\mathcal{L}(X, Y)_n$. A map $\alpha \in \mathfrak{g}$ is then a Lie series if and only if (2.4) holds [see Ref. 30 for an alternative characterization of maps $\alpha : \mathcal{T} \to \mathbb{R}$ that actually belong to $\overline{\mathcal{L}(X, Y)}$].

In particular, the BCH series $Z = Z_1 + Z_2 + Z_3 + \cdots$ given by (1.8) [for X and Y defined as in (2.3)] is a Lie series. From (1.8), it follows that $Z(\bullet) = Z(\circ) = 1$, and for $n = 2, 3, 4, \dots$

TABLE I. First elements E_i of the basis of P. Hall, their corresponding bicolored rooted trees u_i, the values $|i|$, i'', i', $\sigma(u_i)$, and the coefficients $z_i = Z(u_i)/\sigma(u_i)$ in the BCH series (1.9).

i	$|i|$	i'	i''	E_i	u_i	$\sigma(u_i)$	$z_i = \dfrac{Z(u_i)}{\sigma(u_i)}$
1	1	1	0	X	•	1	1
2	1	2	0	Y	○	1	1
3	2	2	1	$[Y,X]$		1	$-\frac{1}{2}$
4	3	3	1	$[[Y,X],X]$		2	$\frac{1}{12}$
5	3	3	2	$[[Y,X],Y]$		1	$-\frac{1}{12}$
6	4	4	1	$[[[Y,X],X],X]$		6	0
7	4	4	2	$[[[Y,X],X],Y]$		2	$\frac{1}{24}$
8	4	5	2	$[[[Y,X],Y],Y]$		2	0
9	5	6	1	$[[[[Y,X],X],X],X]$		24	$-\frac{1}{720}$
10	5	6	2	$[[[[Y,X],X],X],Y]$		6	$-\frac{1}{180}$
11	5	7	2	$[[[[Y,X],X],Y],Y]$		4	$\frac{1}{180}$
12	5	8	2	$[[[[Y,X],Y],Y],Y]$		6	$\frac{1}{720}$
13	5	4	3	$[[[Y,X],X],[X,Y]]$		2	$-\frac{1}{120}$
14	5	5	3	$[[[Y,X],Y],[X,Y]]$		1	$-\frac{1}{360}$

$$nZ(u) = \frac{1}{2}[X - Y, Z](u) + \sum_{p=1}^{[(n-1)/2]} \frac{B_{2p}}{(2p)!}(\mathrm{ad}_Z^{2p}(X+Y))(u) \qquad (2.6)$$

for each $u \in \mathcal{T}$ with $n = |u|$. Recall that, for arbitrary $\alpha, \beta \in \mathfrak{g}$ and $u \in \mathcal{T}$, the value $[\alpha, \beta](u)$ is defined in terms of bicolored rooted trees $u_{(j)}, u^{(j)}$ with less vertices than u, so that (2.6) effectively allows us to compute the values $Z(u)$ for all bicolored rooted trees with arbitrarily high number $|u|$ of vertices. In this way, the characterization (2.4) of maps $\alpha \in \mathfrak{g}$ that are Lie series directly gives a way to write $Z \in \overline{\mathcal{L}(x, Y)}$ in the form (1.9) with

$$z_i = \frac{Z(u_i)}{\sigma(u_i)} \quad \text{for } i \geq 1. \qquad (2.7)$$

For instance, we have according to Table I that in the Hall basis,

$$
\begin{aligned}
Z &= \sum_{i \geq 1} z_i E_i = \sum_{i \geq 1} \frac{Z(u_i)}{\sigma(u_i)} E_i \\
&= Z(\bullet)X + Z(\circ)Y + Z(\text{\scriptsize tree})[Y,X] \\
&\quad + \frac{Z(\text{\scriptsize tree})}{2}[[Y,X],X] + Z(\text{\scriptsize tree})[[Y,X],Y] + \cdots,
\end{aligned}
$$

where the first five coefficients $Z(u_i)$ can be obtained by applying (2.6) with (2.2),

$$[X - Y, Z](\begin{smallmatrix}\bullet\\\circ\end{smallmatrix}) = -Z(\bullet) - Z(\circ) = -2,$$

$$2Z(\begin{smallmatrix}\bullet\\\circ\end{smallmatrix}) = \frac{1}{2}[X - Y, Z](\begin{smallmatrix}\bullet\\\circ\end{smallmatrix}) = -1,$$

$$[X - Y, Z](\begin{smallmatrix}\circ\\\circ\end{smallmatrix}) = 0, \quad 2Z(\begin{smallmatrix}\circ\\\circ\end{smallmatrix}) = 0,$$

$$[X - Y, Z](\bullet\!\!\vee\!\!\bullet) = -2Z(\begin{smallmatrix}\bullet\\\circ\end{smallmatrix}) = 1,$$

$$[X - Y, Z](\bullet\!\!\vee\!\!\circ) = Z(\begin{smallmatrix}\bullet\\\circ\end{smallmatrix}) - Z(\begin{smallmatrix}\circ\\\circ\end{smallmatrix}) = -\frac{1}{2},$$

$$[Z, [Z, X + Y]](\bullet\!\!\vee\!\!\bullet) = -2Z(\bullet)(Z(\bullet) - Z(\circ)) = 0,$$

$$[Z, [Z, X + Y]](\bullet\!\!\vee\!\!\circ) = -Z(\bullet)(Z(\bullet) - Z(\circ)) = 0,$$

$$3Z(\bullet\!\!\vee\!\!\bullet) = \frac{1}{2},$$

$$3Z(\bullet\!\!\vee\!\!\circ) = -\frac{1}{4}. \tag{2.8}$$

In summary, the idea of the formalism is to construct algorithmically a sequence of labeled rooted trees in a one-to-one correspondence with a Hall basis, verifying in addition (2.4). In this way it is quite straightforward to build and characterize Lie series, and, in particular, the BCH series.

B. Detailed treatment

In this subsection we provide a detailed treatment of the main steps involved in the procedure previously sketched, first by analyzing the representation (2.4) of Lie series for the classical Hall basis and then by considering Hall–Viennot bases.

We start by providing an algorithm that constructs the table of values (i', i'') (for $i \geqslant 3$) in (1.10) (together with $|i|$ for $i \geqslant 1$) that determines a classical Hall basis. The algorithm starts by setting

$$1' = 1, \quad 1'' = 0, \quad 2' = 2, \quad 2'' = 0, \quad |1| = 1, \quad |2| = 1,$$
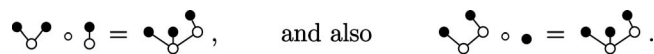
and initializing the counter i as $i = 3$. Then, the values $i', i'', |i|$ for subsequent values of i are set as follows (i^{++} indicates that the value of the counter i is incremented by 1),

Algorithm 1:

> *for* $n = 2, 3, \ldots$
> $\quad j = 1, \ldots, i-1$
> $\quad\quad k = j+1, \ldots, i-1$
> $\quad\quad\quad$ *If* $|j| + |k| = n$ *and* $j \geqslant k''$ *then*
> $\quad\quad\quad\quad i'' = j, \ i' = k, \ |i| = n,$
> $\quad\quad\quad\quad i^{++}.$

The values of $i', i'', |i|$ thus determined satisfy that $i' > i'' \geqslant (i')''$ for $i \geqslant 3$. In addition, $j < i$ if $|j| < |i|$, which implies that $i', i'' < i$ for all $i \geqslant 3$. The values for $|i|$, i', and i'' and the element E_i of the basis for the values of the index i of degree $|i| \leqslant 5$ are displayed in Table I.

On the other hand, it is possible to design a simple recursive procedure to define the bicolored rooted trees u_i appearing in (2.4) in terms of the values of i' and i'' by using the following binary operation. Given $u, v \in \mathcal{T}$, the new rooted tree $u \circ v \in \mathcal{T}$ is a rooted tree with $|u| + |v|$ vertices obtained by grafting the rooted tree v to the root of u (that is to say, $u \circ v$ is a new bicolored rooted tree with the colored vertices of u and v, one edge that makes the root of v a child of the root of u added to the edges of u and v). For instance,

$$\bullet\!\!\vee\!\!\bullet \circ \begin{smallmatrix}\bullet\\\circ\end{smallmatrix} = \bullet\!\!\vee\!\!\bullet^{\bullet}_{\circ}, \quad \text{and also} \quad \bullet\!\!\vee\!\!\circ \circ \bullet = \bullet\!\!\vee\!\!\bullet^{\bullet}_{\circ}.$$

We now define

$$u_1 = \bullet, \quad u_2 = \bigcirc, \quad \text{and} \quad u_i = u_{i'} \circ u_{i''} \quad \text{for } i \geq 3. \tag{2.9}$$

Finally, the symmetry numbers $\sigma_i = \sigma(u_i)$ can also be determined recursively,

$$\sigma_1 = \sigma_2 = 1 \quad \text{and} \quad \sigma_i = \kappa_i \sigma_{i'} \sigma_{i''} \quad \text{for } i \geq 3, \tag{2.10}$$

where $\kappa_i = 1$ if $(i')'' \neq i''$ and $\kappa_i = \kappa_{i'} + 1$ if $(i')'' = i''$.

The bicolored rooted trees u_i, their symmetry numbers $\sigma(u_i)$, and the coefficients $z_i = Z(u_i)/\sigma(u_i)$ in the BCH series (1.9) are displayed in Table I for the first values of the index i, whereas in Table III given in Appendix, the terms of the BCH series (1.9) up to terms of degree of 9 are given in compact form for the classical Hall basis by displaying the values of i', i'', and $z_i = Z(u_i)/\sigma(u_i)$ for each index i.

This procedure can be extended indeed to Hall–Viennot bases. A set $\{E_i : i = 1, 2, 3, \ldots\} \subset \mathcal{L}(X, Y)$ recursively defined as (1.10) with some positive integers i', $i'' < i$ ($i = 3, 4, \ldots$) is a Hall–Viennot basis if there exists a total order relation $>$ in the set of indices $\{1, 2, 3, \ldots\}$ such that $i > i''$ for all $i \geq 3$, and the map

$$d : \{3, 4, \ldots\} \rightarrow \{(j, k) \in \mathbb{Z}^+ \times \mathbb{Z}^+ : j > k \geq j''\}, \tag{2.11}$$

$$d(i) = (i', i'') \tag{2.12}$$

(with the convention $1'' = 2'' = 0$) is bijective.

In Refs. 48 and 36, Hall–Viennot bases are indexed by a subset of words (a Hall set of words) on the alphabet $\{x, y\}$. Such Hall set of words $\{w_i : i \geq 1\}$ can be obtained by defining recursively w_i as the concatenation $w_{i'} w_{i''}$ of the words $w_{i'}$ and $w_{i''}$, with $w_1 = x$ and $w_2 = y$. For instance, the Hall set of words w_i associated with the indices $i = 1, 2, \ldots, 14$ in Table I are x, y, yx, yxx, yxy, $yxxx$, $yxxy$, $yxyy$, $yxxxx$, $yxxxy$, $yxxyy$, $yxyyy$, $yxxyx$, and $yxyyx$.

For the classical Hall basis we have considered before, the map (2.11) has been constructed in such a way that the total order relation $>$ is the natural order relation in $\mathbb{Z}^+$, i.e., $>$ (notice that in Ref. 7 the total order is chosen as $<$).

This is not possible, however, for the Lyndon basis. The Lyndon basis can be constructed as a Hall–Viennot basis by considering the order relation $>$ as follows: $i > j$ if, in lexicographical order (i.e., the order used when ordering words in the dictionary), the Hall word w_i associated with i comes before than the Hall word w_j associated with j. The Hall set of words $\{w_i : i \geq 1\}$ corresponding to the Lyndon basis is the set of Lyndon words, which can be defined as the set of words w on the alphabet $\{x, y\}$ satisfying that, for arbitrary decompositions of w as the concatenation $w = uv$ of two nonempty words u and v, the word w is smaller than v in lexicographical order.[48,25]

Now, the representation (2.4) of a map $\alpha \in \overline{\mathcal{L}(X, Y)}$ [and, in particular, the BCH series (1.9) with (2.7)] for any Hall–Viennot basis can be stated as follows.

Theorem 2.1: *Given a total order relation $>$ in $\mathbb{Z}^+$ and a bijection (2.11) satisfying that $i > i''$ for all $i \geq 3$, then any map $\alpha \in \overline{\mathcal{L}(X, Y)}$ admits the representation (2.4) for the Hall basis (1.10) and the bicolored rooted trees u_i and their symmetry numbers $\sigma_i = \sigma(u_i)$ recursively defined as (2.9) and (2.10).*

Theorem 2.1 can be proven as a corollary of Theorem 3 and Remark 17 in Ref. 30. Actually, in Ref. 30 it is shown that (2.4) holds for a different set $\hat{\mathcal{T}} = \{u_1, u_2, u_3, \ldots\}$ of bicolored rooted trees associated with a Hall basis, for which $\sigma(u_i) = 1$ for all i. However, the set of Hall rooted trees we consider here (which is the set $\hat{\mathcal{T}}^*$ considered in Remark 17 in Ref. 30) has some advantages from the computational point of view.

In Table II, we display the elements E_i of the Lyndon basis with degree $|i| \leq 5$, the corresponding Lyndon words w_i, the bicolored rooted trees u_i, the values $|i|$, i'', i', $\sigma(u_i)$, and the coefficients $z_i = Z(u_i)/\sigma(u_i)$ in the BCH series (1.9).

TABLE II. First elements E_i of the Lyndon basis, their corresponding Lyndon words w_i and bicolored rooted trees u_i, the values $|i|$, i'', i', $\sigma(u_i)$, and the coefficients $z_i = Z(u_i)/\sigma(u_i)$ in the BCH series (1.9).

i	$|i|$	i'	i''	w_i	E_i	u_i	$\sigma(u_i)$	$z_i = \frac{Z(u_i)}{\sigma(u_i)}$
1	1	1	0	x	X		1	1
2	1	2	0	y	Y		1	1
3	2	1	2	xy	$[X,Y]$		1	$\frac{1}{2}$
4	3	3	2	xyy	$[[X,Y],Y]$		2	$\frac{1}{12}$
5	3	1	3	xxy	$[X,[X,Y]]$		1	$\frac{1}{12}$
6	4	4	2	$xyyy$	$[[[X,Y],Y],Y]$		6	0
7	4	1	4	$xxyy$	$[X,[[X,Y],Y]]$		2	$\frac{1}{24}$
8	4	1	5	$xxxy$	$[X,[X,[X,Y]]]$		1	0
9	5	6	2	$xyyyy$	$[[[[X,Y],Y],Y],Y]$		24	$\frac{1}{720}$
10	5	5	3	$xxyxy$	$[[X,[X,Y]],[X,Y]]$		2	$\frac{1}{360}$
11	5	3	4	$xyxyy$	$[[X,Y],[[X,Y],Y]]$		2	$\frac{1}{120}$
12	5	1	6	$xxyyy$	$[X,[[[X,Y],Y],Y]]$		6	$\frac{1}{180}$
13	5	1	7	$xxxyy$	$[X,[X,[[X,Y],Y]]]$		2	$\frac{1}{180}$
14	5	1	8	$xxxxy$	$[X,[X,[X,[X,Y]]]]$		1	$-\frac{1}{720}$

C. Practical aspects in the implementation

An important ingredient in the whole procedure is the practical implementation of the Lie bracket $[\alpha,\beta]$ of two Lie series $\alpha,\beta \in \overline{\mathcal{L}(X,Y)} \subset \mathfrak{g}$, which we address next. Let us consider for each $u \in \mathcal{T}$ the sequence

$$S(u) = \{(u_{(1)},u^{(1)}), \ldots, (u_{(|u|-1)},u^{(|u|-1)})\} \tag{2.13}$$

of pairs of bicolored rooted trees used to define the Lie bracket $[\alpha,\beta]$ in (2.1). For instance,

$$S(\ ) = \left\{(\ ,\ \circ),(\bullet,\ ),(\bullet,\ )\right\}.$$

It can be seen that the sequences $S(u)$ satisfy the following recursion. If $u = v \circ w$, where $v,w \in \mathcal{T}$, then, let $p = |v|-1$, $q = |w|-1$, and

$$S(v) = \{(v_{(1)},v^{(1)}), \ldots, (v_{(p)},v^{(p)})\}, \quad S(w) = \{(w_{(1)},w^{(1)}), \ldots, (w_{(q)},w^{(q)})\},$$

then

$$S(u) = \{(w,v),(v_{(1)},v^{(1)} \circ w), \ldots, (v_{(p)},v^{(p)} \circ w),(w_{(1)},v \circ w^{(1)}), \ldots, (w_{(q)},v \circ w^{(q)})\}. \tag{2.14}$$

From the point of view of implementation, it is important to observe that, if one wants to compute the BCH formula (1.9) up to terms of degree $|i| \leq n$, there is no need to compute all values $Z(u)$ for $u \in T$ with $|u| \leq n$, since only the values (2.7) of $Z(u)$ for the rooted trees u_i associated with the elements E_i of the basis are used in (1.9). However, when applying the recursion (2.6) for the bicolored rooted trees u_i, $i = 1,2,3,\ldots$ (that we will call Hall rooted trees), one finds that $Z(u)$ needs to be computed for some additional bicolored rooted trees u. Observe, for instance, that in (2.8), the value $Z(u)$ for $u = $ was needed in order to get the value $Z($ $)$ from the recursion (2.6). This is due to the fact that in (2.1), the bicolored rooted trees $u_{(j)}$ need not be Hall rooted trees when u is a Hall rooted tree ($u^{(j)}$ is, however, necessarily a Hall rooted tree in that case).

The minimal set $\widetilde{\mathcal{T}}_n$ of bicolored rooted trees u for which $Z(u)$ needs to be computed in order to get the values of $Z(u_i)$ for Hall rooted trees with $|i| \leqslant n$ by using recursion (2.6) can be determined by requiring that

$$\{u_i : |i| \leqslant n\} \subset \widetilde{\mathcal{T}}_n \subset \mathcal{T} \quad \text{and} \quad S(\widetilde{\mathcal{T}}_n) \subset \widetilde{\mathcal{T}}_n \times \widetilde{\mathcal{T}}_n.$$

It can be seen that the subset $\widetilde{\mathcal{T}}_n$ of bicolored rooted trees can be alternatively defined as follows: We say that a bicolored rooted tree $v \in \mathcal{T}$ is covered by $u \in \mathcal{T}$ if either v can be obtained from u by removing some of its vertices and edges or $u = v$. For instance, the bicolored rooted trees covered by the tree u_{11} in Table II are

$$\bullet \bullet, \quad \circ, \quad \overset{\bullet}{\bullet}, \quad \overset{\circ}{\bullet}, \quad \vee\!\!\bullet, \quad \vee\!\!\circ, \quad \overset{\bullet}{\underset{\bullet}{<}}, \quad \overset{\circ}{\underset{\bullet}{<}}, \quad \overset{\bullet}{\underset{\bullet}{<}}\!\!\bullet, \quad \overset{\circ}{\underset{\bullet}{<}}\!\!\circ.$$

Then, it can be seen that $\widetilde{\mathcal{T}}_n$ is the set of bicolored rooted trees covered by some of the trees in $\{u_i : |i| \leqslant n\}$.

As a summary of this treatment, we next describe the main steps of the algorithm that we use to compute the BCH series up to terms of a given degree N for an arbitrary Hall–Viennot basis. Let m_N be sum of the dimensions of the homogeneous subspaces $\mathcal{L}(X,Y)_n$ for $1 \leqslant n \leqslant N$ and let $\widetilde{m}_N$ be the number of bicolored rooted trees in $\widetilde{\mathcal{T}}_N$ (so that $m_N \leqslant \widetilde{m}_N$). We proceed as follows for a given N:

(1) Determine the values i', i'' for each $i = 1, \ldots, m_N$ such that the E_i given by (1.10) are the elements of degree $|i| \leqslant N$ of the required Hall–Viennot basis. Algorithm 1 can be used in the case of the basis of P. Hall. We use a similar (although slightly more complex) algorithm for the general case.

(2) Determine the bicolored rooted trees $u \in \widetilde{T}_N$ together with the $|u| - 1$ pairs of bicolored rooted trees in $S(u)$ recursively obtained by (2.14). Actually, we associate each bicolored rooted tree in $\widetilde{\mathcal{T}}_N$ with a positive integer, such that $\widetilde{\mathcal{T}}_N = \{u_i : i = 1, 2, \ldots, \widetilde{m}_N\}$ (and $\{u_i : i = 1, 2, \ldots, m_N\}$ is the set of Hall trees of degree $|i| \leqslant N$). Each $S(u_i)$ is then represented as a list of $|i| - 1$ pairs of positive integers.

(3) Represent the truncated versions of Lie series α (truncated up to terms of degree N) as a list of $\widetilde{m}_N$ real values $(\alpha_1, \ldots, \alpha_{\widetilde{m}_N})$ corresponding to $(\alpha(u_1), \ldots, \alpha(u_{\widetilde{m}_N}))$. The Lie bracket $\gamma = [\alpha, \beta]$ of two Lie series can be implemented as a way to obtain the list $(\gamma_1, \ldots, \gamma_{\widetilde{m}_N})$ from the lists $(\alpha_1, \ldots, \alpha_{\widetilde{m}_N})$ and $(\beta_1, \ldots, \beta_{\widetilde{m}_N})$ in terms of the pairs of integers representing $S(u_i)$ for each $i = 1, \ldots, \widetilde{m}_N$.

(4) Represent the truncated versions of BCH series Z (truncated up to terms of degree N) as a list of $\widetilde{m}_N$ rational values $(Z_1, \ldots, Z_{\widetilde{m}_N})$ corresponding to $(Z(u_1), \ldots, Z(u_{\widetilde{m}_N}))$, which can be obtained by initializing that list as $(1, 1, 0, \ldots, 0)$ and applying (2.6) repeatedly for $n = 2, \ldots, N$.

It is worth noticing that the number of trees in $\widetilde{\mathcal{T}}_n$ is different for different Hall–Viennot bases. For instance, for the basis of P. Hall, $\widetilde{T}_{20}$ has 724 018 bicolored rooted trees, while for the Lyndon basis the set $\widetilde{T}_{20}$ has 1 952 325 bicolored rooted trees. Due to this fact, the amount of memory and CPU time required to compute with our algorithm the BCH formula up to a given degree for the Lyndon basis is considerably larger than for the basis of P. Hall. Moreover, the number of nonzero coefficients z_i in the BCH formula differs considerably in both bases. For instance, there are 109 697 nonvanishing coefficients z_i (out of 111 013 elements E_i of degree $|i| \leqslant 20$) in the BCH formula for the basis of P. Hall, while for the Lyndon basis the number of nonvanishing coefficients z_i is 76 760.

III. OPTIMAL CONVERGENCE DOMAIN OF THE BCH SERIES

A. The BCH formula and the Magnus expansion

One particularly simple way of obtaining a sharp bound on the convergence domain for the BCH series consists in relating it with the Magnus expansion for linear differential equations. For the sake of completeness, we summarize here the main features of this procedure.

Suppose we have the nonautonomous linear differential equation

$$\frac{dU}{dt} = A(t)U, \quad U(0) = I, \tag{3.1}$$

where $U(t)$ and $A(t)$ are operators acting on some Hilbert space $\mathcal{H}$ (in particular, $n \times n$ real or complex matrices). Then the idea is to express the solution $U(t)$ as the exponential of a certain operator $\Omega(t)$,

$$U(t) = \exp \Omega(t). \tag{3.2}$$

By substituting (3.2) into (3.1), one can derive the differential equation satisfied by the exponent Ω,

$$\Omega' = \sum_{k=0}^{\infty} \frac{B_k}{k!} \mathrm{ad}_{\Omega}^k(A(t)), \quad \Omega(0) = O. \tag{3.3}$$

By applying Picard's iteration on (3.3), one gets an infinite series for $\Omega(t)$,

$$\Omega(t) = \sum_{m=1}^{\infty} \Omega_m(t), \tag{3.4}$$

whose terms can be obtained recursively from

$$\Omega_1(t) = \int_0^t A(t_1)dt_1,$$

$$\Omega_m(t) = \sum_{j=1}^{m-1} \frac{B_j}{j!} \int_0^t (\mathrm{ad}_{\Omega(s)}A(s))_m ds, \quad m \geq 2. \tag{3.5}$$

Equations (3.2) and (3.4) constitute the so-called Magnus expansion for the solution of (3.1), whereas the infinite series (3.4) with (3.5) is known as the *Magnus series*.

Since the 1960s,[49] the Magnus expansion has been successfully applied as a perturbative tool in numerous areas of physics and chemistry, from atomic and molecular physics to nuclear magnetic resonance and quantum electrodynamics (see Refs. 4 and 5 for a review and a list of references). Also, since the work by Iserles and Nørsett,[20] it has been used as a tool to construct practical algorithms for the numerical integration of Eq. (3.1), while preserving the main qualitative properties of the exact solution.

In general, the Magnus series does not converge unless A is small in a suitable sense, and several bounds to the actual radius of convergence have been obtained along the years. Recently, the following theorem has been proven.[9]

Theorem 3.1: *Let us consider the differential equation $U' = A(t)U$ defined in a Hilbert space $\mathcal{H}$, dim $\mathcal{H} < \infty$, with $U(0) = I$, and let $A(t)$ be a bounded linear operator on $\mathcal{H}$. Then, the Magnus series $\Omega(t) = \Sigma_{k=1}^{\infty}\Omega_k(t)$, with Ω_k given by (3.5) converges in the interval $t \in [0,T)$ such that*

$$\int_0^T \|A(s)\|ds < \pi$$

and the sum $\Omega(t)$ *satisfies* $\exp \Omega(t) = U(t)$. *The statement also holds when* $\mathcal{H}$ *is infinite dimensional if* U *is a normal operator (in particular, if* U *is unitary). Here* $\|\cdot\|$ *stands for the norm defined by the inner product on* $\mathcal{H}$.

Moreover, it has been shown that the convergence domain of the Magnus series provided by this theorem is the best result one can get for a generic bounded operator $A(t)$ in a Hilbert space, in the sense that it is possible to find specific $A(t)$ where the series diverges for any time t such that $\int_0^t \|A(s)\| ds > \pi$.[29,9]

Now, given two operators X and Y, let us consider Eq. (3.1) with

$$A(t) = \begin{cases} Y, & 0 \le t < 1 \\ X, & 1 \le t \le 2. \end{cases} \tag{3.6}$$

Clearly, the exact solution at $t=2$ is given by $U(2) = e^X e^Y$. On the other hand, if we apply recurrence (3.5) to compute $U(2)$ with the Magnus expansion, $U(2) = e^{\Omega(2)}$, we get $\Omega_1(2) = X + Y$ and more generally $\Omega_n(2) = Z_n$ in (1.4). In other words, the BCH series can be considered as the Magnus expansion corresponding to the differential equation (3.1) with $A(t)$ given by (3.6) at $t = 2$.

Since $\int_0^{t=2} \|A(s)\| ds = \|X\| + \|Y\|$, Theorem 3.1 leads to the following bound on the convergence of the BCH series.

Theorem 3.2: *Let X and Y be two bounded elements in a Hilbert space $\mathcal{H}$ with $\dim \mathcal{H} \ge 2$. Then the BCH formula in the form (1.4), i.e., expressed as a series of homogeneous Lie polynomials in X and Y, converges when $\|X\| + \|Y\| < \pi$.*

Of course, this result can be generalized to any set $X_1, X_2, \ldots, X_k$ of bounded operators: the corresponding BCH series is convergent if $\|X_1\| + \cdots + \|X_k\| < \pi$ in the 2-norm.

Let us illustrate the result provided by Theorem 3.2 with a simple example involving 2×2 matrices.

Example 1: Given

$$X = \begin{pmatrix} \alpha & 0 \\ 0 & -\alpha \end{pmatrix}, \quad Y = \begin{pmatrix} 0 & \beta \\ 0 & 0 \end{pmatrix}, \tag{3.7}$$

with $\alpha, \beta \in \mathbb{C}$, a simple calculation shows that

$$\log(e^X e^Y) = X + \frac{2\alpha}{1 - e^{-2\alpha}} Y,$$

which is an analytic function for $|\alpha| < \pi$ with first singularities at $\alpha = \pm i\pi$. Therefore, the BCH formula cannot converge if $|\alpha| \ge \pi$, independently of $\beta \ne 0$. By taking the spectral norm, it is clear that $\|X\| = |\alpha|$, $\|Y\| = |\beta|$, so that the convergence domain given by Theorem 3.2 is $|\alpha| + |\beta| < \pi$. Notice that in the limit $|\beta| \to 0$ this domain is optimal. $\qquad \square$

Generally speaking, however, the bound given by Theorem 3.2 is conservative, i.e., the BCH series converges for larger values of $\|X\|$ and $\|Y\|$. Thus, in the previous example, for any α and β with $|\alpha| < \pi$ and $|\alpha| + |\beta| \ge \pi$, the BCH series also converges. One would like therefore to have a more realistic characterization of this feature. It turns out that this is indeed feasible for complex $n \times n$ matrices.

B. Convergence for matrices

1. Convergence determined by the eigenvalues

For complex $n \times n$ matrices it is possible to use the theory of analytic matrix functions and more specifically, the logarithm of an analytic matrix function, in a similar way as in the Magnus expansion,[9] to characterize more precisely the convergence of the BCH series.

To begin with, let us introduce a parameter $\varepsilon \in \mathbb{C}$ and consider the substitution $(X, Y) \mapsto (\varepsilon X, \varepsilon Y)$ into Eq. (1.1). It is clear that

$$U(\varepsilon) \equiv e^{\varepsilon X} e^{\varepsilon Y}$$

is an analytic function of ε, $\det U(\varepsilon) \neq 0$ and the matrix function $Z(\varepsilon) = \log U(\varepsilon)$ is also analytic at $\varepsilon = 0$. Equivalently, the series $Z(\varepsilon)$ is convergent for sufficiently small ε. It turns out that the actual radius of convergence of this series is related with the existence of multiple eigenvalues of $U(\varepsilon)$. Let us denote by $\rho_1(\varepsilon), \ldots, \rho_n(\varepsilon)$ the eigenvalues of the matrix $U(\varepsilon)$. Observe that $U(0) = I$, so that $\rho_1(0) = \cdots = \rho_n(0) = 1$, and we can take the principal values of the logarithm, $\log \rho_1(0) = \cdots = \log \rho_n(0) = 0$. In essence, if the analytic matrix function $U(\varepsilon)$ has an eigenvalue $\rho_0(\varepsilon_0)$ of multiplicity $l > 1$ for a certain ε_0 such that (a) there is a curve in the ε-plane joining $\varepsilon = 0$ with $\varepsilon = \varepsilon_0$, and (b) the number of equal terms in $\log \rho_1(\varepsilon_0)$, $\log \rho_2(\varepsilon_0), \ldots, \log \rho_l(\varepsilon_0)$ such that $\rho_k(\varepsilon_0) = \rho_0$, $k = 1, \ldots, l$ is less than the maximum dimension of the elementary Jordan block corresponding to ρ_0, then the radius of convergence of the series $Z(\varepsilon) = \Sigma_{k \geqslant 1} \varepsilon^k Z_k$ verifying $\exp Z(\varepsilon) = U(\varepsilon)$ is precisely $r = |\varepsilon_0|$.[9]

More specifically, we find first the values of the parameter ε for which the characteristic polynomial $\det(U(\varepsilon) - \rho I)$ has multiple roots and write them in order of nondecreasing absolute value,

$$\varepsilon_0^{(1)}, \varepsilon_0^{(2)}, \varepsilon_0^{(3)}, \ldots . \tag{3.8}$$

Next, we consider the circle $|\varepsilon| = |\varepsilon_0^{(1)}|$ in the complex ε-plane and denote by $\rho_0^{(1)}$ an eigenvalue of $U(\varepsilon_0^{(1)})$ with multiplicity $l_1 > 1$. Let ε move along some fixed curve L from $\varepsilon = 0$ to $\varepsilon = \varepsilon_0^{(1)}$ in the circle $|\varepsilon| \leqslant |\varepsilon_0^{(1)}|$. Then it is clear that l_1 eigenvalues $\rho_j(\varepsilon)$ will tend to $\rho_0^{(1)}$ at $\varepsilon = \varepsilon_0^{(1)}$. If these points lie at $\varepsilon = \varepsilon_0^{(1)}$ on the same sheet of the Riemann surface of the function $\log z$, and this is true for all (possible) multiple eigenvalues of $U(\varepsilon)$ at $\varepsilon = \varepsilon_0^{(1)}$, then $\varepsilon_0^{(1)}$ is called a *extraneous root*. Otherwise, $\varepsilon_0^{(1)}$ is called a *nonextraneous root*.

By the analysis carried out in Ref. 51, when $|\varepsilon| < |\varepsilon_0^{(1)}|$ the numbers $\log \rho_j(\varepsilon)$ are uniquely determined as eigenvalues of the matrix $Z(\varepsilon)$ and this series is convergent. This is also true at $|\varepsilon| = |\varepsilon_0^{(1)}|$ if $\varepsilon_0^{(1)}$ is an extraneous root, since then the eigenvalues of $Z(\varepsilon)$ retain their identity throughout the collision process, so that we proceed to the next value in the sequence (3.8) until a nonextraneous root is obtained.

Assume, for simplicity, that $\varepsilon_0^{(2)}$ is the first nonextraneous root, for which there exists an eigenvalue ρ_0 of $U(\varepsilon)$ with multiplicity $l > 1$. Associated with this multiple eigenvalue ρ_0 there is a pair of integers (p, q) defined as follows.

The integer p is the greatest number of equal terms in the set of numbers $\log \rho_1(\varepsilon_0)$, $\log \rho_2(\varepsilon_0), \ldots, \log \rho_l(\varepsilon_0)$ such that $\rho_k(\varepsilon_0) = \rho_0$, $k = 1, \ldots, l$.

The integer q is the maximum degree of the elementary divisors $(\rho - \rho_0)^k$ of $U(\varepsilon_0)$, i.e., the maximum dimension of the elementary Jordan block corresponding to ρ_0.

Under these conditions, it has been proven that if $p < q$ for the eigenvalue ρ_0, then the radius of convergence of the series $Z(\varepsilon) = \Sigma_{k \geqslant 1} \varepsilon^k Z_k$ is precisely $r = |\varepsilon_0|$.[51]

Although in some cases with $p \geqslant q$ the series $Z(\varepsilon)$ may converge at $|\varepsilon| = |\varepsilon_0|$ and the radius of convergence r is greater than $|\varepsilon_0|$ (for instance, when X and Y are diagonal), this situation is exceptional in a topological sense, as explained in Ref. 51, pp. 65 and 66.

2. Examples

In order to illustrate this result we next consider a pair of examples involving also 2×2 matrices.

Example 1: The first example involves again the matrices X and Y given by (3.7). In this case

$$U(\varepsilon) = e^{\varepsilon X} e^{\varepsilon Y} = \begin{pmatrix} e^{\varepsilon \alpha} & \varepsilon \beta e^{\varepsilon \alpha} \\ 0 & e^{-\varepsilon \alpha} \end{pmatrix}.$$

The first values of ε for which there are multiple eigenvalues of $U(\varepsilon)$ are

$$\varepsilon = 0, \quad \varepsilon = \pm\, i\frac{\pi}{\alpha}.$$

The first value, $\varepsilon = 0$, is clearly an extraneous root, whereas the eigenvalues of the matrix $U(\varepsilon)$ move along the unit circle, one clockwise and the other counterclockwise from

$$\rho_{1,2}(0) = 1 \text{ to } \rho_{1,2}(i\pi/\alpha) = -1$$

when ε varies along the imaginary axis from $\varepsilon = 0$ to $\varepsilon = i\pi/\alpha$ (the same considerations apply to the case $\varepsilon = -i\pi/\alpha$). Then, obviously, $p = 1$ and $q = 2$, so that the radius of convergence of the series $Z(\varepsilon)$ is

$$|\varepsilon| = \frac{\pi}{|\alpha|}.$$

By fixing $\varepsilon = 1$, we get the actual domain of convergence of the BCH series as $|\alpha| = \pi$, i.e., the same result as in Sec. III A   □

Example 2: Consider now the matrices

$$A = \begin{pmatrix} 0 & 0 \\ 1 & 0 \end{pmatrix}, \quad B = \begin{pmatrix} 0 & 1 \\ 0 & 0 \end{pmatrix},$$

and $X = \alpha A$, $Y = \alpha B$, with $\alpha > 0$. Then

$$U(\varepsilon) = \begin{pmatrix} 1 & \alpha\varepsilon \\ \alpha\varepsilon & 1 + \alpha^2\varepsilon^2 \end{pmatrix} \tag{3.9}$$

has multiple eigenvalues when $\varepsilon_0^{(1)} = 0$, $\varepsilon_0^{(2)} = \pm\, i2/\alpha$. As ε varies along the imaginary axis from $\varepsilon = 0$ to $\varepsilon = \varepsilon_0^{(2)}$, the eigenvalues of the matrix $U(\varepsilon)$,

$$\rho_{1,2}(\varepsilon) = 1 + \frac{\alpha^2}{2}\varepsilon^2 \pm \sqrt{\left(1 + \frac{\alpha^2}{2}\varepsilon^2\right)^2 - 1},$$

move along the unit circle, one clockwise and the other counterclockwise from

$$\rho_{1,2}(0) = 1 \text{ to } \rho_{1,2}(\varepsilon_0^{(2)}) = -1.$$

Thus, $\rho_1(\varepsilon_0^{(2)})$ and $\rho_2(\varepsilon_0^{(2)})$ lie on different sheets of the Riemann surface of the function $\log z$ and therefore $\varepsilon_0^{(2)}$ is a nonextraneous root, with $p = 1$. Since $U(\varepsilon_0^{(2)}) \neq -I$, we have $q = 2$, so that the radius of convergence of the series $Z(\varepsilon)$ is precisely

$$r = |\varepsilon_0^{(2)}| = \frac{2}{\alpha}. \tag{3.10}$$

This result should be compared with the bound provided by the Magnus expansion. Since $\|A\| = \|B\| = 1$, Theorem 3.2 guarantees the convergence of the BCH series in this case whenever $2\alpha|\varepsilon| < \pi$ or $|\varepsilon| < \pi/2\alpha$, which, in view of (3.10), is clearly a conservative estimate.

We can also check numerically the rate of convergence of the BCH series in this example as a function of the parameter ε. Let us denote by $Z^{[N]}$ the sum of the first N terms of the series, i.e.,

$$Z^{[N]}(\varepsilon) = \sum_{n=1}^{N} Z_n(\varepsilon)$$

and compute, for $\alpha = 2$ and different values of ε, the matrix

$$E_r(\varepsilon) = U(\varepsilon)e^{-Z^{[N]}(\varepsilon)} - I,$$

where $U(\varepsilon)$ is given by (3.9). If ε belongs to the convergence domain of the BCH series for the matrices X and Y (i.e., $|\varepsilon| < 1$), then $E_r(\varepsilon) \rightarrow 0$ as $N \rightarrow \infty$.

First we take $\varepsilon = \frac{1}{4}$. With $N=10$, the elements of E_r are of order of 10^{-7}, whereas adding five additional terms in the series, $N=15$, the elements of E_r are approximately 10^{-10}.

Next we choose $\varepsilon = 0.9$, i.e., a value near the boundary of the convergence domain. In this case with $N=15$ the convergence of the series does not manifest at all. In fact, a much larger number of terms is required to achieve significant results. Thus, for the elements of E_r to be of order of 10^{-8} we need to compute $N=150$ terms of the BCH series, whereas with $N=200$ the elements of E_r are of order of 10^{-10}. The computations have been carried out with the recurrence (1.8). $\square$

As this example clearly shows, it is not always possible to determine accurately the convergence domain of the BCH series by computing successive approximations, since the rate of convergence can be slow indeed near the boundary. For this reason it could be of interest to design a procedure to apply in practice the characterization of the convergence in terms of the eigenvalues of the matrix $U(\varepsilon)$ analyzed in Sec. III B 1 for matrices.

This procedure could be as follows. Given two matrices X, Y, take the product of exponentials

$$U(\varepsilon) = e^{\varepsilon X} e^{\varepsilon Y}$$

with $\varepsilon = re^{i\theta}$. Next, define a grid in the ε-plane, for instance, in polar coordinates (r, θ), by $\Delta r = r_f/(n+1)$, $\Delta\theta = 2\pi/(m+1)$ for two integers n, $m \geq 1$ and a sufficiently large value $r_f > 1$. Then, for each point in the grid $(r_k = k\Delta r, \ \theta_l = l\Delta\theta)$, $k = 1, \ldots, n+1$, $l = 0, 1, \ldots, m$, compute the corresponding matrix $U(\varepsilon)$ and its eigenstructure, locating where there are multiple eigenvalues (within a prescribed tolerance). If some of these multiple eigenvalues have a negative real part, there exists a point in the neighborhood where the conditions enumerated in Sec. III B 1 are satisfied, and therefore we have approximately located the value of ε where the BCH series fails to converge. This approximation can be made more accurate by applying, for instance, Newton's method. The actual radius of convergence will be given by the smallest number r found in this way. Finally, if $r > 1$, then obviously the BCH series corresponding to X and Y converges.

IV. SOME APPLICATIONS

As an illustration of the usefulness of the previous results, in this section we present two not so trivial applications of the formalism developed in Sec. II for constructing explicitly the BCH series up to arbitrarily high order.

A. The symmetric BCH formula

Sometimes it is necessary to compute the Lie series W defined by

$$\exp\left(\tfrac{1}{2}X\right)\exp(Y)\exp\left(\tfrac{1}{2}X\right) = \exp(W). \tag{4.1}$$

This occurs, for instance, if one is interested in obtaining the order conditions satisfied by time-symmetric composition methods for the numerical integration of differential equations.[52,39] Two applications of the usual BCH formula give then the expression of W in the Hall basis of $\mathcal{L}(X, Y)$.

A more efficient procedure is obtained, however, by introducing a parameter τ in (4.1) such that

$$W(\tau) = \log(e^{\tau X/2} e^Y e^{\tau X/2}) \tag{4.2}$$

and deriving the differential equation satisfied by $W(\tau)$. From the derivative of the exponential map, one gets

$$\frac{dW}{d\tau} = X + \sum_{n=2}^{\infty} \frac{B_n}{n!}\mathrm{ad}_W^n X, \quad W(0) = Y,$$ (4.3)

whence it is possible to construct explicitly W as the series $W(\tau) = \Sigma_{k=0}^{\infty} W_k(\tau)$, with

$$W_1(\tau) = X\tau + Y,$$

$$W_2(\tau) = 0,$$

$$W_l(\tau) = \sum_{j=2}^{l-1} \frac{B_j}{j!} \int_0^\tau (\mathrm{ad}_W^j X)_l ds, \quad l \geq 3,$$ (4.4)

where, in general, $W_{2m} = 0$ for $m \geq 1$. By following a similar approach as with Eq. (1.8) in the usual BCH series in Sec. II, the recursion (4.4) allows one to express W in (4.1) as

$$W = \sum_{i \geq 1} w_i E_i.$$ (4.5)

The coefficients w_i of this series up to degree of 9 in the classical Hall basis are collected in Table IV in Appendix. As with the usual BCH series, the coefficients up to degree of 19 in both Hall and Lyndon bases can be found at www.gicas.uji.es/research/bch.html.

With respect to the convergence of the series, theorem (3.2) guarantees that W is convergent at least when $\|X\| + \|Y\| < \pi$.

B. The BCH formula and a problem of Thompson

In a series of papers,[43,32,44,45] Thompson considered the problem of constructing a representation of the BCH formula as

$$e^X e^Y = e^Z, \quad \text{with} \quad Z = SXS^{-1} + TYT^{-1},$$ (4.6)

for certain functions $S = S(X,Y)$ and $T = T(X,Y)$ depending on X and Y. By using analytic techniques related with the Kashiwara–Vergne method, Rouvière[38] proved that a Lie series $\rho(X,Y)$ exists such that

$$S = e^{\rho(X,Y)}, \quad T = e^{\rho(-Y,-X)}$$ (4.7)

and converges when X, Y are replaced by normed elements near 0, whereas the representation (4.6) is global when both X and Y are skew-Hermitian matrices.[43]

Thompson himself developed a computational technique for constructing explicitly the series $\rho(X,Y)$ up to terms of degree of 10. Although his results were not published, he pointed out that they furnished strong evidence of the convergence of the series $\rho(X,Y)$ on the closed unit sphere in any norm for which $\|[X,Y]\| \leq \|X\|\|Y\|$.[45]

With the aim of clarifying this issue and illustrating the techniques developed in Sec. II, we proceed next to compute $\rho(X,Y)$. Since $\rho(X,Y) \in \overline{\mathcal{L}(X,Y)}$, i.e., is a Lie series, it can be written as

$$\rho(X,Y) = \sum_{i \geq 1} \rho_i E_i,$$

where the elements E_i have been introduced in (1.9), and the goal is to determine the coefficients ρ_i. This can be accomplished as follows. From the well known formula $e^U V e^{-U} = e^{\mathrm{ad}_U} V$, it is clear that

$$Z = e^{\mathrm{ad}_{\rho(X,Y)}}X + e^{\mathrm{ad}_{\rho(-Y,-X)}}Y. \tag{4.8}$$

Next we expand $e^{\mathrm{ad}_{\rho(X,Y)}}X$ and $e^{\mathrm{ad}_{\rho(-Y,-X)}}Y$ into infinite series as a linear combination of the Hall basis in $\mathcal{L}(X,Y)$ and match the resulting terms with the corresponding to the BCH series for Z. Then a recursive system of equations is obtained for the coefficients ρ_i.

It is, in fact, possible to get a closed expression for $\rho(X,Y)$ up to terms Y^2 by taking into account the corresponding formula of Z.[34] Specifically, from

$$Z = X + \frac{\mathrm{ad}_X}{1 - e^{-\mathrm{ad}_X}}Y \bmod Y^2, \tag{4.9}$$

a simple calculation leads to

$$\rho(X,Y) = f(\mathrm{ad}_X)Y \bmod Y^2,$$

with the function $f(z)$ given by

$$f(z) = \frac{e^z}{1 - e^z} + \frac{1}{z}e^{z/4} = -\frac{1}{4} - \frac{5}{96}z + \frac{1}{384}z^2 + \frac{143}{92\,160}z^3 + \frac{1}{122\,880}z^4 + \cdots . \tag{4.10}$$

Working in the classical Hall basis, the complete expression up to degree of 4 reads

$$f(z) = -\tfrac{1}{4}Y + \tfrac{5}{96}[Y,X] + \tfrac{1}{384}[[Y,X],X] + \tfrac{11}{768}[[Y,X],Y] - \tfrac{143}{92\,160}[[[Y,X],X],X] - \tfrac{283}{92\,160}[[[Y,X],X],Y]$$
$$+ \tfrac{11}{23\,040}[[[Y,X],Y],Y],$$

i.e., the corresponding equations have a unique solution. This is not the case, however, at degree of 5, where a free parameter appears, which can be chosen to be ρ_{10}. Then

$$\rho_{12} = \frac{-137 - 184\,320\rho_{10}}{184\,320}, \quad \rho_{13} = \frac{-511 - 737\,280\rho_{10}}{737\,280}.$$

As a matter of fact, if higher degrees are considered, more and more free parameters appear in the corresponding solution. Thus, at degree of 7 there are two additional parameters (for instance, ρ_{26} and ρ_{30}), whereas at degree of 8 ρ_{50} and ρ_{52} can be chosen as free parameters. We conclude, therefore, that there are infinite solutions to the problem posed by Thompson depending on an increasing number of free parameters. An interesting issue would be to determine the value of these parameters in order to render the whole series convergent on a domain as large as possible.

C. Distribution of coefficients in the Lyndon basis

As we previously mentioned, there are noteworthy differences in the results obtained when the algorithm of Sec. II is applied to the BCH series in the classical Hall basis and the Lyndon basis, particularly with respect to the number of vanishing coefficients. In the basis of P. Hall there are 1316 zero coefficients out of 111 013 up to degree $m=20$, whereas in the Lyndon basis the number of vanishing terms rises to 34 253 (more than 30% of the total number of coefficients).

More remarkably, one notices that the distribution of these vanishing coefficients in the Lyndon basis follows a very specific pattern. Before entering into the details, let us denote for simplicity $\mathcal{L}_m \equiv \mathcal{L}(X,Y)_m$. We first remark that, for each $m \geq 2$, the Lyndon basis $\mathcal{B}_m$ of $\mathcal{L}_m$ is a disjoint union $\mathcal{B}_m = \mathcal{B}_{m,1} \cup \mathcal{B}_{m,2}$ with $\mathcal{B}_{m,2} = [X, \mathcal{B}_{m-1}]$. Thus, $\mathcal{L}_m = \mathcal{L}_{m,1} \oplus \mathcal{L}_{m,2}$, where $\mathcal{L}_{m,2} = [X, \mathcal{L}_{m-1}]$, and $\mathcal{B}_{m,k}$ $(k=1,2)$ is a basis of $\mathcal{L}_{m,k}$. In particular, $\mathrm{ad}_X^{m-1} Y \in \mathcal{B}_m$. In this sense, from our computations we make two observations. First, the coefficient in the BCH formula of the element $\mathrm{ad}_X^{m-1} Y$ in the basis $\mathcal{B}_m$ is 0 for even m. Second, the coefficients for the terms in $\mathcal{B}_{m,1}$ are also zero for even m. This gives a total number of

TABLE III. Table of values of i' and i'' for $i \geq 3$ in (1.10) for the classical Hall basis and the values $z_i \in \mathbb{Q}$ in the BCH formula (1.9).

i	i'	i''	z_i	i	i'	i''	z_i	i	i'	i''	z_i
1	1	0	1	44	25	2	1/10 080	87	31	3	−11/30 240
2	2	0	1	45	26	2	23/120 960	88	32	3	−19/100 800
3	2	1	−1/2	46	27	2	1/10 080	89	33	3	−1/43 200
4	3	1	1/12	47	28	2	1/60 480	90	34	3	−1/10 080
5	3	2	−1/12	48	29	2	0	91	35	3	−1/50 400
6	4	1	0	49	15	3	0	92	15	4	−1/33 600
7	4	2	1/24	50	16	3	1/40 032	93	16	4	−13/120 960
8	5	2	0	51	17	3	23/30 240	94	17	4	−1/10 080
9	6	1	−1/720	52	18	3	1/2 240	95	18	4	−11/201 600
10	6	2	−1/180	53	19	3	1/15 120	96	19	4	−1/43 200
11	7	2	1/180	54	20	3	0	97	20	4	−1/7 560
12	8	2	1/720	55	21	3	1/2 250	98	21	4	−1/10 080
13	4	3	−1/120	56	22	3	1/10 080	99	22	4	1/50 400
14	5	3	−1/360	57	9	4	0	100	23	4	1/20 160
15	9	1	0	58	10	4	1/10 080	101	15	5	−23/302 400
16	9	2	−1/1 440	59	11	4	−1/20 160	102	16	5	−1/5 760
17	10	2	−1/360	60	12	4	−1/20 160	103	17	5	13/151 200
18	11	2	−1/1 440	61	13	4	0	104	18	5	19/120 960
19	12	2	0	62	14	4	−1/2 520	105	19	5	1/33 600
20	6	3	0	63	9	5	1/4 032	106	20	5	−13/30 240
21	7	3	−1/240	64	10	5	1/840	107	21	5	−23/100 800
22	8	3	−1/720	65	11	5	1/1 440	108	22	5	−1/100 800
23	5	4	1/240	66	12	5	1/12 096	109	23	5	−1/33 600
24	15	1	1/30 240	67	13	5	1/1 260	110	9	6	−1/60 480
25	15	2	1/5 040	68	14	5	1/10 080	111	10	6	−1/90 720
26	16	2	1/3 780	69	7	6	−1/10 080	112	11	6	1/30 240
27	17	2	−1/3 780	70	8	6	−13/30 240	113	12	6	−11/302 400
28	18	2	−1/5 040	71	8	7	−1/3 360	114	13	6	1/15 120
29	19	2	−1/30 240	72	42	1	−1/1 209 600	115	14	6	1/3 780
30	9	3	1/2 016	73	42	2	−1/151 200	116	9	7	−11/120 960
31	10	3	23/15 120	74	43	2	−1/56 700	117	10	7	−1/6 720
32	11	3	1/5 040	75	44	2	−1/75 600	118	11	7	−1/14 400
33	12	3	−1/10 080	76	45	2	1/75 600	119	12	7	−11/120 960
34	13	3	1/1 260	77	46	2	1/56 700	120	13	7	−1/20 160
35	14	3	1/5 040	78	47	2	1/151 200	121	14	7	17/100 800
36	6	4	1/5 040	79	48	2	1/1 209 600	122	9	8	−1/20 160
37	7	4	−1/10 080	80	24	3	−1/43 200	123	10	8	17/151 200
38	8	4	1/1 680	81	25	3	−37/302 400	124	11	8	1/6 048
39	6	5	13/15 120	82	26	3	−11/60 480	125	12	8	1/60 480
40	7	5	−1/1 120	83	27	3	−11/302 400	126	13	8	−1/100 800
41	8	5	−1/5 040	84	28	3	11/302 400	127	14	8	1/37 800
42	24	1	0	85	29	3	1/100 800				
43	24	2	1/60 480	86	30	3	−1/7 560				

$$n_c(2p) = \dim(\mathcal{L}_{2p}) - \dim(\mathcal{L}_{2p-1}) + 1, \quad p \geq 2,$$

vanishing coefficients of terms of degree $m = 2p$ in the BCH formula written in the Lyndon basis. Thus, for instance, when $p = 10$, the number of total number of vanishing coefficients is $n_c(20) = \dim(\mathcal{L}_{20}) - \dim(\mathcal{L}_{19}) + 1 = 52\,377 - 27\,594 + 1 = 24\,784$.

TABLE IV. Table of values of i' and i'' for $i \geq 3$ in (1.10) for the classical Hall basis and the values $w_i \in \mathbb{Q}$ in the symmetric BCH formula (4.1)

i	i'	i''	w_i	i	i'	i''	w_i	i	i'	i''	w_i
1	1	0	1	44	25	2	0	87	31	3	1/4 608
2	2	0	1	45	26	2	0	88	32	3	23/134 400
3	2	1	0	46	27	2	0	89	33	3	1/37 800
4	3	1	−1/24	47	28	2	0	90	34	3	1/23 040
5	3	2	−1/12	48	29	2	0	91	35	3	1/201 600
6	4	1	0	49	15	3	0	92	15	4	193/6 451 200
7	4	2	0	50	16	3	0	93	16	4	53/483 840
8	5	2	0	51	17	3	0	94	17	4	25/193 536
9	6	1	7/5 760	52	18	3	0	95	18	4	1/22 400
10	6	2	7/1 440	53	19	3	0	96	19	4	−13/1 209 600
11	7	2	1/180	54	20	3	0	97	20	4	53/483 840
12	8	2	1/720	55	21	3	0	98	21	4	17/161 280
13	4	3	1/480	56	22	3	0	99	22	4	−3/44 800
14	5	3	−1/360	57	9	4	0	100	23	4	−19/322 560
15	9	1	0	58	10	4	0	101	15	5	367/4 838 400
16	9	2	0	59	11	4	0	102	16	5	193/645 120
17	10	2	0	60	12	4	0	103	17	5	247/604 800
18	11	2	0	61	13	4	0	104	18	5	53/241 920
19	12	2	0	62	14	4	0	105	19	5	1/33 600
20	6	3	0	63	9	5	0	106	20	5	53/161 280
21	7	3	0	64	10	5	0	107	21	5	193/403 200
22	8	3	0	65	11	5	0	108	22	5	13/201 600
23	5	4	0	66	12	5	0	109	23	5	−1/5 600
24	15	1	−31/967 680	67	13	5	0	110	9	6	11/774 114
25	15	2	−31/161 280	68	14	5	0	111	10	6	1/290 304
26	16	2	−13/30 240	69	7	6	0	112	11	6	−1/15 360
27	17	2	−53/120 960	70	8	6	0	113	12	6	−89/1 209 600
28	18	2	−1/5 040	71	8	7	0	114	13	6	−11/241 920
29	19	2	−1/30 240	72	42	1	127/154 828 800	115	14	6	−13/80 640
30	9	3	−53/161 280	73	42	2	127/19 353 600	116	9	7	1/12 096
31	10	3	−11/12 096	74	43	2	157/7 257 600	117	10	7	11/64 512
32	11	3	−3/4 480	75	44	2	367/9 676 800	118	11	7	1/33 600
33	12	3	−1/10 080	76	45	2	23/604 800	119	12	7	−11/120 960
34	13	3	−1/4 032	77	46	2	79/3 628 800	120	13	7	1/35 840
35	14	3	−1/6 720	78	47	2	1/151 200	121	14	7	−29/134 400
36	6	4	−19/80 640	79	48	2	1/1 209 600	122	9	8	211/1 935 360
37	7	4	−1/10 080	80	24	3	367/19 353 600	123	10	8	173/604 800
38	8	4	17/40 320	81	25	3	473/4 838 400	124	11	8	5/24 192
39	6	5	−53/60 480	82	26	3	41/215 040	125	12	8	1/60 480
40	7	5	−19/13 440	83	27	3	211/1 209 600	126	13	8	61/403 200
41	8	5	−1/5 040	84	28	3	89/1 209 600	127	14	8	−1/151 200
42	24	1	0	85	29	3	1/100 800				
43	24	2	0	86	30	3	79/967 680				

With these considerations in mind, we can proceed next to explain the observed phenomena. First, notice that expression (4.9) gives explicitly the last term of the BCH series in the Lyndon basis at each degree. By formally expanding in power series of ad_X we get

$$Z = X + Y + \frac{1}{2}\text{ad}_X Y + \sum_{k=2}^{\infty} (-1)^k \frac{B_k}{k!}\text{ad}_X^k Y \bmod Y^2.$$

Since $B_{2n+1}=0$ for all $n \geq 1$, the coefficient of $\text{ad}_X^k Y$ is nonvanishing only for even values of k, or equivalently, for odd values of the degree m.

As for the remaining zero coefficients, let us consider at this point the symmetric BCH formula (4.1) again. Clearly the series (4.5) only contains terms of odd degree, i.e., $W = \Sigma_{i \geq 0} W_{2i+1}$, where $W_i \in \mathcal{L}_i$. By denoting $P = X/2$ and forming the composition $\exp(P)\exp(W)\exp(-P)$ one gets trivially

$$e^P e^W e^{-P} = e^X e^Y = e^Z,$$

i.e., the standard BCH formula. In the terminology of dynamical systems, $\exp(W)$ and $\exp(Z)$ are said to be conjugated. Alternatively, we can write $\exp(Z)=\exp(\text{ad}_P)\exp(W)$, so that $Z = \exp(\text{ad}_P)W$. It is worth to write explicitly this relation for each term $Z_m \in \mathcal{L}_m$ of the series $Z = \Sigma_{m \geq 0} Z_m$ by separating the odd and even degree cases. Specifically,

$$Z_{2p+1} = W_{2p+1} + \sum_{j=1}^{p} \frac{1}{(2j)!)2^{2j}}\text{ad}_X^{2j} W_{2p-2j+1},$$

$$Z_{2p} = \sum_{j=1}^{p} \frac{1}{(2j-1)!)2^{2j-1}}\text{ad}_X^{2j-1} W_{2p-2j+1}.$$

From these expressions, it is clear that Z_{2p+1} contains terms in the whole subspace $\mathcal{L}_{2p+1,1} \oplus \mathcal{L}_{2p+1,2}$ (due to the presence of W_{2p+1}), whereas Z_{2p} belongs to the subspace $\mathcal{L}_{2p,2}$, whose dimension is equal to $\dim(\mathcal{L}_{2p-1})$. In other words, the remaining $\dim(\mathcal{L}_{2p}) - \dim(\mathcal{L}_{2p-1})$ *must* necessarily vanish. In this sense, the Lyndon basis seems the natural choice to get systematically the BCH series with the minimum number of terms. Nevertheless, compared to the basis of P. Hall, more CPU time and memory are required to compute the BCH with our algorithm in the Lyndon basis. In particular, 1.5 Gbytes are required to compute the BCH formula up to degree of 20 in the Hall basis, whereas 3.6 Gbytes of memory are needed in the Lyndon basis.

V. CONCLUDING REMARKS

The effective computation of the BCH series has a long history and is closely related with the more general problem of carrying out symbolic computations in free Lie algebras. In this work we have presented a new algorithm which allows us to get a closed expression of the series $Z = \log(e^X e^Y)$ up to degree of 20 in terms of an arbitrary Hall–Viennot basis of the free Lie algebra generated by X and Y, $\mathcal{L}(X,Y)$, requiring reasonable computational resources. As far as we know, no other results are available up to this degree in terms of a basis of $\mathcal{L}(X,Y)$. The algorithm is based on some more general results presented in Ref. 30 on the connection of labeled rooted trees with an arbitrary Hall–Viennot basis of the free Lie algebra.

We have carried out explicitly the computations to get the coefficients of the BCH series in terms of both the classical Hall basis and the Lyndon basis, with some noteworthy differences in the corresponding results, as analyzed in Sec. IV C.

We have also addressed the problem of the convergence of the series when X and Y are replaced by normed elements. In the particular case of X and Y being matrices, we have provided a characterization of the convergence in terms of the eigenvalues of e^Z.

Although here we have considered only the BCH series, it is clear that other more involved calculations can be done, as is illustrated, for instance, by the problem of Thompson studied in Sec. IV B. As a matter of fact, we intend to develop a general purpose package to carry out symbolic computations in a free Lie algebra generated by more than two operators.

ACKNOWLEDGMENTS

APPENDIX: COEFFICIENTS OF THE BCH FORMULA

In Table III we collect the indices i' and i'' for $i \geq 3$ in (1.10) for the classical Hall basis and the values of the coefficients z_i in the BCH formula (1.9) up to degree of 9, whereas in Table IV we gather the corresponding coefficients for the symmetric BCH formula (4.1).

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