The non-commutative Korteweg–de Vries hierarchy and combinatorial Pöppe algebra

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Abstract We give a constructive proof, to all orders, that each member of the non-commutative potential Korteweg–de Vries hierarchy is a Fredholm Grassmannian flow and is therefore linearisable. Indeed we prove this for any linear combination of fields from this hierarchy. That each member of the hierarchy is linearisable, and integrable in this sense, means that the time evolving solution can be generated from the solution to the corresponding linear dispersion equation in the hierarchy, combined with solving an associated linear Fredholm equation representing the Marchenko equation. Further, we show that within the class of polynomial partial differential fields, at every order, each member of the non-commutative potential Korteweg-de Vries hierarchy is unique. Indeed, we prove to all orders, that each such member matches the non-commutative Lax hierarchy field, which is therefore a polynomial partial differential field. We achieve this by constructing the abstract combinatorial algebra that underlies the non-commutative potential Korteweg-de Vries hierarchy. This algebra is the non-commutative polynomial algebra over the real line generated by the set of all compositions endowed with the Pöppe product. This product is the abstract representation of the product rule for Hankel operators pioneered by Ch. Pöppe for integrable equations such as the Sine-Gordon and Korteweg-de Vries equations. Integrability of the hierarchy members translates, in the combinatorial algebra, to proving the existence of a ‘Pöppe polynomial’ expansion for basic compositions in terms of ‘linear signature expansions’. Proving the existence of such Pöppe polynomial expansions boils down to solving a linear algebraic problem for the expansion coefficients, which we solve constructively to all orders.

Keywords KdV hierarchy · Grassmannian flow · combinatorial Pöppe algebra

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* The author dedicates this paper to Claudia Wulff, a wonderful mathematician, person and brave friend, who passed away on the 12th of June 2021.
1 Introduction

We prove, to all orders, that each member of the non-commutative potential Korteweg–de Vries hierarchy is a Fredholm Grassmannian flow. As such, each member of the hierarchy is linearisable, and therefore integrable in this sense. The proof straightforwardly extends to any linear combination of fields from this hierarchy. In addition we show that within the class of non-commutative polynomial partial differential fields, each member of the hierarchy is unique. Indeed, we prove that each such member matches the non-commutative Lax hierarchy field, which at every order is therefore a non-commutative polynomial partial differential field. By ‘field’ here we mean the right-hand side in the partial differential equation when we express it in the form of the derivative in time of the solution equals the linear and nonlinear terms on the right. By a ‘non-commutative polynomial partial differential field’ we mean that the field can expressed as a non-commutative multivariate polynomial whose arguments are the solution and spatial partial derivatives of the solution only. When we say each member of the hierarchy is linearisable, we mean that the time-evolving solution can be generated by solving the corresponding linear dispersion equation in the hierarchy for the ‘scattering data’, and the solution to the hierarchy is then generated by solving an associated linear Fredholm equation involving the scattering data that represents a Marchenko equation. Ch. Pöppe pioneered this approach to finding solutions to classical scalar integrable systems such as the Sine-Gordon equation, Korteweg–de Vries equation and its scalar hierarchy, as well as the nonlinear Schrödinger equation in a sequence of papers. See Pöppe [55, 56, 57], Pöppe and Sattinger [58] and Bauhardt and Pöppe [8]. Recently Doikou et al. [22] streamlined Pöppe’s approach, demonstrating that only Pöppe’s celebrated kernel product rule is required to establish the linearisation of the classical Korteweg–de Vries and nonlinear Schrödinger equations. Subsequently Doikou et al. [21] demonstrated the approach, as considered by Bauhardt and Pöppe [8], naturally extends to the non-commutative nonlinear Schrödinger and Korteweg–de Vries equations. Malham [46] then used this approach to establish the linearisation of the non-commutative fourth order quintic nonlinear Schrödinger equation. The results herein extend Pöppe’s approach to the whole non-commutative Korteweg–de Vries hierarchy. However the algebraic structures we develop to establish the linearisation of this hierarchy, provide a deeper insight into the ‘integrability’ properties of the hierarchy, not only establishing uniqueness, but also transporting the question of integrability into the existence of polynomial expansions in a given combinatorial algebra equipped with an abstract version of the Pöppe product. In honour of Ch. Pöppe, we call this combinatorial algebra the Pöppe algebra and the polynomial expansions, Pöppe polynomials.

Let us now explain the approach in some more detail. Consider a non-commutative nonlinear partial differential equation for $g = g(x, t)$ of the form:

$$
\partial_t g = d(\partial)g + \pi(g, \partial g, \partial^2 g, \ldots),
$$

where $\partial = \partial_x$. Here the ‘field’ which is the right-hand side consists of the linear $d(\partial)g$ terms, where $d(\partial)$ is a constant coefficient odd homogeneous polynomial with the lowest degree term of degree 3, and the nonlinear terms $\pi$. We suppose $\pi$ is a precise homogeneous non-commutative polynomial function of the arguments indicated up to including derivatives of $g$ of order $\deg(d) - 2$. We assume we have absorbed the linear terms in $\pi$ into $d(\partial)g$, so the lowest degree term in $\pi$ is 2. Now suppose the matrix valued
function \( p = p(x, t) \) satisfies the corresponding linear partial differential equation:

\[ \partial_t p = d(\partial)p. \]

The quantity \( p = p(x, t) \) represents the so-called ‘scattering data’. Pöppe elevated the Marchenko equation to the operator level, associating a Hilbert–Schmidt Hankel operator \( P = P(x, t) \) with \( p \) as follows. For any square integrable function \( \phi \) set:

\[ (P\phi)(y; x, t) := \int_{-\infty}^{0} p(y + z + x; t)\phi(z)\,dz. \]

The operator \( P = P(x, t) \) satisfies the linear operator differential equation,

\[ \partial_t P = d(\partial)P, \]

which we augment with the following linear Fredholm equation for \( G = G(x, t) \),

\[ P = G(id - P). \]

This represents the Marchenko equation at the operator level. These are the only ingredients we need for the non-commutative potential Korteweg–de Vries hierarchy.

Pöppe’s kernel product rule involving Hankel operators is then key. Suppose \( F = F(x, t) \) and \( F' = F'(x, t) \) are both Hilbert–Schmidt operators that are continuously dependent on the parameters \( x \) and \( t \). Further suppose that \( H = H(x, t) \) and \( H' = H'(x, t) \) are Hilbert–Schmidt Hankel operators, that are continuously differentially dependent on \( x \) and \( t \). Let \( [F] \) denote the kernel of any Hilbert–Schmidt operator \( F \). Then the Fundamental Theorem of Calculus implies Pöppe’s kernel product rule:

\[ [F\partial_x(HH')F'](y; x, t) = [FH](y, 0; x, t)[H'H'](0; z; x, t). \]

With this in hand, we just compute \( \partial_t G - d(\partial)G \), where \( G \) satisfies the linear Fredholm equation above, and apply the kernel bracket operator \([\cdot]\). If we write \( G = PU \) where \( U := (id - P)^{-1} \), then we see that basic calculus properties such as \( \partial U = U(\partial P)U \) generate nonlinear terms. The goal is then to use only the kernel product rule to establish ‘closed forms’ for the nonlinear terms generated. By this we mean that the terms generated can be expressed as a constant coefficient non-commutative polynomial in \( \partial[G], \partial^2[G] \) and so forth. Note we do not include \([G]\) itself in this polynomial as our aim is to derive the non-commutative potential Korteweg–de Vries hierarchy equations.

Fleshing this out further, since \( U := (id - P)^{-1} \), we know \( U \equiv id + UP \). We have seen that \( \partial U = U(\partial P)U \). The general Leibniz rule implies:

\[ \partial^n U = U(\partial^n P)U + n(\partial U)(\partial^{n-1} P)U + \cdots + n(\partial^{n-1} U)(\partial P)U, \]

where we have combined the \( \partial^n U \) terms appearing in the Leibniz expansion, and used the definition for \( U \). Suppose with \( n \geq 3 \) odd, and we take \( d(\partial) = \mu_n \partial^n P \) with \( \mu_n \) a real constant. Then since \( G = PU \) and \( U \equiv id + UP \equiv id + PU \), we observe \( \partial_t G = \partial_t U = U(\partial_t P)U = \mu_n U(\partial^n P)U \) and so,

\[ \partial_t G = \mu_n(\partial^n U) - n(\partial U)(\partial^{n-1} P)U - \cdots - n(\partial^{n-1} U)(\partial P)U \]

\[ \Rightarrow \partial_t[G] = \mu_n\partial^n[G] - \mu_n\left(n(\partial U)(\partial^{n-1} P)U) + \cdots + n((\partial^{n-1} U)(\partial P)U)ight), \]

where we have used that \( \partial^n G = \partial^n U \) and so \( \partial^n[G] = [\partial^n U] \). The question is, can we express all the remaining terms on the right shown in the desired ‘closed form’? In
other words, can we express them as a constant coefficient non-commutative polynomial in \( \partial[G], \partial^2[G] \) and so forth. Note that \( \mu_n [U(\partial^n P)U] = \partial_1[G] \) and from the general Leibniz rule above, we have,

\[
[U(\partial^n P)U] = \partial^n[U] - n [(\partial U)(\partial^{n-1}P)U] - \cdots - n [(\partial^{n-1}U)(\partial P)U].
\]

Thus an equivalent question is, can we express all the terms on the right in terms of a constant coefficient non-commutative polynomial in \( \partial[U], \partial^2[U] \) and so forth? This is obviously trivial in the case of the first term on the right.

Let us consider a quick example, naturally for the case \( n = 3 \). It is possible to show in a few lines, see Example 2 in Section 3 that

\[
[U(\partial^3 P)U] = \partial^3[U] - 3 (\partial[U]) (\partial[U]).
\]

Note, from the Pöppe kernel product rule above, when we write \( \partial[U] \) we mean \( (\partial[U])(y,0;x,t) (\partial[U])(0,z;x,t) \). This thus establishes that \( [G] = [G](y, z; x, t) \) satisfies the partial differential equation,

\[
\partial_1[G] = \mu_3 \partial^3[G] - 3 \mu_3 (\partial[G]) (\partial[G]).
\]

If we set \( y = z = 0 \), then as a special case, \( g(x,t) := [G](0,0;x,t) \) satisfies

\[
\partial g = \mu_3 \partial^3 g - 3 \mu_3 (\partial g)(\partial g),
\]

which is the non-commutative potential Korteweg-de Vries equation.

It is the equivalent question above that is crucial to tackling the higher order cases systematically. Returning to the original general Leibniz rule for operators above, we observe the expression for \( \partial^n U \) contains all the lower order derivatives \( \partial^{n-1}U, \ldots, \partial U \) as factors in the terms on the right. We can insert the corresponding Leibniz rules for the lower order terms. This generates the expansion, see Theorems 1 and 2

\[
\partial^n U = \sum \chi(a_1a_2\cdots a_k) UP_{a_1}UP_{a_2}U \cdots UP_{a_k}U,
\]

where the sum is over all compositions \( a_1a_2\cdots a_k \in C(n) \) of \( n \), and \( P_k := \partial^k P \). We call the real coefficient \( \chi(a_1a_2\cdots a_k) \) the signature character, it has an explicit form as a product of Leibniz coefficients; see Definition 5.

Applying the kernel bracket to this expression we naturally get,

\[
[\partial^n U] = [UP_n U] + \sum \chi(a_1a_2\cdots a_k) [UP_{a_1}UP_{a_2}U \cdots UP_{a_k}U],
\]

where the sum is over all compositions \( a_1a_2\cdots a_k \in C(n) \backslash \{n\} \); we exclude \( n \) itself which we have distinguished in the first term and used that \( \chi(n) = 1 \). Further, in this context, Pöppe’s kernel product rule is equivalent to the following statement. If \( u = a_1 \cdots a_k \) and \( v = b_1 \cdots b_l \) are compositions and \( a \) and \( b \) are natural numbers, then denoting \( m_u := UP_{a_1}UP_{a_2}U \cdots UP_{a_k}U \) and \( m_v := UP_{b_1}UP_{b_2}U \cdots UP_{b_l}U \), we have

\[
[m_u P_{a}U][P_{b}m_v] = [m_u P_{a+1}UP_{b}m_v] + [m_u P_{a}UP_{b+1}m_v] + 2 [m_u P_{a}UP_{b}P_{1}U]m_v.
\]

Now let us take stock. We see that we can always express the terms \( [\partial^n U] \) for any \( n \in \mathbb{N} \) as a linear combination of monomials of the form \( [UP_{a_1}UP_{a_2}U \cdots UP_{a_k}U] \). Further we have an explicit expression for the product of any such monomials. Indeed, we observe that in principle we can focus our attention entirely on an algebra of such
monomials in which the product is the one just stated. The terms \([\partial^n U]\) are simple, natural example linear combinations of monomials in this algebra. We call them the \textit{linear signature expansions}, or \textit{signature expansions} for short. Now, rearranging the expression for \([\partial^n U]\) above and noting that \(\chi(n) = 1\), we find,

\[
[UP_nU] = [\partial^n U] - \sum \chi(a_1a_2 \cdots a_k) [UP_{a_1}UUP_{a_2}U \cdots UP_{a_k}U],
\]

where the sum is over all compositions \(a_1a_2 \cdots a_k \in \mathcal{C}(n)\}\{\{n\}\}. And now the \textit{equivalent question} we asked previously above becomes, can we express \([UP_nU]\), or equivalently all the terms on the right, in terms of a constant coefficient non-commutative polynomial in \([\partial U]\), \([\partial^2 U]\), and so forth?

We observe that the variable labels ‘\(P\)’ and ‘\(U\)’ are superfluous. In our discussion above, we ultimately consider the real algebra of the monomials \([UP_{a_1}U \cdots UP_{a_k}U]\) where the product of any two such monomials is given by the Pöppe kernel product. We can replace these monomials directly by compositions \(a_1a_2 \cdots a_k \in \mathcal{C}\), with the corresponding Pöppe product ‘\(*\)’ given as follows. For two compositions \(ua\) and \(bv\) in \(\mathcal{C}\), where we distinguish the last letter and first letters \(a\) and \(b\) as indicated, the Pöppe product for compositions is given by

\[
(ua) * (bv) = u(a + 1)bv + ua(b + 1)v + 2 \cdot (ua1bv).
\]

This mirrors the Pöppe kernel product rule above. We further observe that we can confine ourselves to computations in the algebra \(\mathbb{R}\langle \mathcal{C}\rangle\), the real algebra over \(\mathcal{C}\) endowed with the product ‘\(*\)’. We define the linear signature expansions in this context as the elements \(n \in \mathbb{R}\langle \mathcal{C}\rangle\) given for any \(n \in \mathbb{N}\) by,

\[
n = \sum_{w \in \mathcal{C}(n)} \chi(w) \cdot w.
\]

These correspond to the terms \(\partial^n [U]\). Naturally the product of any two linear signature expansions \(m\) and \(n\) in \(\mathbb{R}\langle \mathcal{C}\rangle\) is given by, \(m * n = \sum \chi(u)\chi(v) \cdot u * v\), where the sum is over all \(u \in \mathcal{C}(m)\) and \(v \in \mathcal{C}(n)\). Further from our coding just above, the single letter composition \(n \in \mathbb{R}\langle \mathcal{C}\rangle\) corresponds to the term \([UP_nU]\). Now the \textit{equivalent question} above becomes: for every odd \(n \in \mathbb{N}\), can we find a polynomial \(\pi_n = \pi_n(1, 2, \ldots, n)\) such that \(\pi_n(1, 2, \ldots, n - 2, n) = n\) where the polynomial \(\pi_n\) has the form:

\[
\pi_n := \frac{1}{2(n+1)} \sum_{k=1}^{n} \sum_{a_1a_2 \cdots a_k \in \mathcal{C}^*(n)} c_{a_1a_2 \cdots a_k} \cdot a_1 \ast a_2 \ast \cdots \ast a_k,
\]

where \(\mathcal{C}^*(n) \subset \mathcal{C}(n)\) represents the subset of compositions \(w = a_1a_2 \cdots a_k\) of \(n\) such that \(a_1 + a_2 + \cdots + a_k = n - k + 1\). This is the appropriate subset of \(k\)-factor monomials of the form \(a_1 \ast a_2 \ast \cdots \ast a_k\) that generate compositions in \(\mathcal{C}(n)\) once their Pöppe products are evaluated. The coefficients \(c_{a_1a_2 \cdots a_k}\) are real constants. These polynomials are the \textit{Pöppe polynomials} mentioned earlier. To determine the coefficients \(c_{a_1a_2 \cdots a_k}\), we expand all the Pöppe products in all the monomials \(a_1 \ast a_2 \ast \cdots \ast a_k\) in \(\pi_n\). Equating the coefficients of all the compositions in \(\mathcal{C}(n)\) thus generated, results in an overdetermined system of linear equations for the coefficients \(c_{a_1a_2 \cdots a_k}\) for all \(a_1a_2 \cdots a_k \in \mathcal{C}^*(n)\). The right-hand side of the equation \(\pi_n = n\) is simply \(n \in \mathbb{R}\langle \mathcal{C}\rangle\) and thus all the linear equations are homogeneous apart from that corresponding to the
single letter composition \( n \in \mathbb{R} \langle C \rangle_* \), i.e. corresponding to the coefficient \( c_n \). We show there is a unique solution to this linear system and thus the corresponding nonlinear partial differential equation in the hierarchy is integrable. In Section 7 at each odd order \( n \in \mathbb{N} \), we derive the non-commutative Lax hierarchy from scratch in the algebra \( \mathbb{R} \langle C \rangle_* \). Using a technique analogous to that just outlined, we show that for every odd \( n \in \mathbb{N} \), the polynomial expansions for the single letter compositions \( n \in \mathbb{R} \langle C \rangle_* \) generated by the non-commutative iterative Lax procedure can also be expressed as a polynomial expansion in terms of \( 1, 2, \ldots, (n - 2) \) and \( n \). Since we previously proved such an expansion for the single letter compositions \( n \in \mathbb{R} \langle C \rangle_* \) is unique, the non-commutative nonlinear partial differential equations generated by the non-commutative Lax hierarchy, and those we generated using the Pöppe polynomials just above, are one and the same. Further, as part of this whole procedure, we also establish a new co-algebra, the ‘signature co-algebra’.

The algebraic approach to the non-commutative Korteweg–de Vries hierarchy introduced herein, necessarily arose from the desire to extend the Grassmannian flow solution approach for classical integrable systems. This approach has its roots in the work of Dyson [23], Ablowitz et al. [1] and Mumford [52], and in particular in a series of papers by Pöppe [53,54,55,56,57], Pöppe and Sattinger [58] and Bauhardt and Pöppe [8]. Other fundamental approaches include the scheme by Zakharov and Shabat [60,61] as well as, for example, the unified transform method, see Fokas and Pelloni [29]. Pöppe’s approach was recently revisited by McKean [50]. Combining the Grassmannian flow solution approach for infinite dimensional Riccati partial differential systems developed in Beck et al. [10,11], with the Hankel operator approach for classical integrable systems developed by Pöppe, Doikou et al. [22] showed how for general initial data, time evolutionary solutions to the Korteweg–de Vries and nonlinear Schrödinger equations represent Grassmannian flows. Such flows are linearisable in the sense outlined above. Stylianidis [62] and Doikou et al. [21] extended this approach to the non-commutative versions of these classical integrable systems. Malham [46] extended the approach to the non-commutative fourth order quintic nonlinear Schrödinger equation. Early work on non-commutative integrable systems includes that by Fordy and Kulisch [30], Ablowitz et al. [3,12] and Ercolani and McKean [26]. Work by Nijhoff et al. [53], Fu [31] and Fu and Nijhoff [22] and the direct linearisation technique they develop is extremely close in spirit to the approach we develop herein. Even closer, in the sense that they elevate the linearisation solution technique to the operator level and seek solutions in the form of inverse Fredholm operators associated with scattering data exactly as we do, is the recent work by Carillo and Schoenlieb [14,15,16]. Also see Aden and Carl [6], Hamanaka and Toda [43], as well as Sooman [65]. The work by Treves [67,68] is also close in the sense of seeking to represent the nonlinear fields in the Korteweg–de Vries hierarchy as polynomials in a given algebra of formal pseudo-differential symbols. Also see Buryak and Rossi [13]. General algebraic solution approaches to the Kadomtsev–Petviashvili hierarchies in the context of weakly nonassociative algebras, which are close to the approach we adopt herein, can be found in Dimakis and Müller–Hoissen [24], whereas Dimakis and Müller–Hoissen [14] consider closer connections to shuffle and Rota–Baxter algebras. See Reutenauer [60], Malham and Wiese [47], Ebrahimi–Fard et al. [24] and also Ebrahimi–Fard et al. [25] for more on shuffle algebras and references therein for Rota-Baxter algebras. For very recent work on the non-commutative Korteweg–de Vries equation, see Degasperis and Lombardo [18], Pelinovsky and Stepanyants [54] and Adamopoulou and Papamikos [5]. Hankel operators have also received a lot of recent attention, see for example Grudsky and Rybkin [25,27], Grellier and Ger-
ard [35], and Blower and Newsham [12]. The connection between Grassmannians and classical integrable systems was first explored by Sato [61, 62] and developed further by Segal and Wilson [63]. Also see Kasman [44, 45] and Hamanaka and Toda [43].

To summarise, what we achieve herein is that, for the non-commutative potential Korteweg-de Vries hierarchy:

(i) We show that establishing integrability at any odd order $n$ is equivalent to determining the existence of certain polynomial expansions for the single letter composition $n$, in the real algebra of compositions equipped with the Pöppe product, in terms of natural so-called (linear) signature expansions of compositions.

Using this algebraic structure, we prove to all orders that each member of the hierarchy:

(ii) Is a Fredholm Grassmannian flow and thus linearisable, and so integrable in this sense. We establish this for any linear combination of fields from the hierarchy;

(iii) Is unique within the class of polynomial partial differential fields;

(iv) Matches the non-commutative Lax hierarchy field, which is therefore a polynomial partial differential field.

We also:

(v) Establish a new co-algebra we call the signature co-algebra based on the de-Pöppe co-product, which tensorially decomposes any composition into the sum of all possible composition pairs that produced that composition via the Pöppe product.

Our paper is organised as follows. In Section 2 we introduce Pöppe’s formulation for integrable system via Hankel operators. We discuss why solutions generated using this approach are example Fredholm Grassmannian flows. We begin our process of abstraction in Section 4, showing how all the inherent operators in Pöppe’s formulation, including the key operator corresponding to the inverse Fredholm operator associated with the scattering data, can be expressed as an expansion in the concatenation algebra of compositions over the real field. After applying the kernel bracket operator to this concatenation algebra, we show how, necessarily, we need to consider the real algebra of compositions equipped with the abstract form of the Pöppe product. A short Section 5 follows, in which we introduce the signature co-algebra. Though strictly not required for the subsequent analysis and proofs in the sections which follow, our motivation and the results therein help elucidate some of the constructs we use and rely on in them. Our main result is stated and proved in Section 6. Indeed, at each odd order we prove the existence of a Pöppe polynomial expansion and thus the integrability of the corresponding non-commutative nonlinear partial differential equation. We generate the non-commutative Lax hierarchy in Section 7. We prove that the non-commutative hierarchy of equations we established in the preceding section, correspond precisely to those in the non-commutative Lax hierarchy. Lastly in Section 8 we explore extensions of the results we establish herein, and in particular, briefly, the application of this approach to the non-commutative nonlinear Schrödinger hierarchy.

2 Pöppe formulation and Grassmannian flows

Let us begin by defining and deriving the main analytical ingredients we need. We consider Hilbert–Schmidt integral operators which depend on both a spatial parameter
$x \in \mathbb{R}$ and a time parameter $t \in [0, \infty)$. In this section $\partial_t$ represents the partial derivative with respect to the time parameter $t$ while $\partial = \partial_x$ represents the partial derivative with respect to the spatial parameter $x$. For any Hilbert–Schmidt operator $F = F(x,t)$ there exists a square-integrable kernel $f = f(y, z; x, t)$ such that for any square-integrable function $\phi$,

$$(F \phi)(y; x, t) = \int_{-\infty}^{\partial} f(y, z; x, t) \phi(z) \, dz.$$

**Definition 1 (Kernel bracket)** With reference to the operator $F$ just above, we use the kernel bracket notation $[F]$ to denote the kernel of $F$:

$$[F](y; z; x, t) := f(y, z; x, t).$$

Hilbert–Schmidt Hankel operators play a critical role herein. We consider Hankel operators which depend on a parameter $x$ as follows.

**Definition 2 (Hankel operator with parameter)** We say a given time-dependent Hilbert–Schmidt operator $H$ with corresponding square-integrable kernel $h$ is Hankel or additive with parameter $x \in \mathbb{R}$ if its action, for any square-integrable function $\phi$, is

$$(H \phi)(y; x, t) := \int_{-\infty}^{\partial} h(y + z + x; t) \phi(z) \, dz.$$

Such Hankel operators are a key facet of Pöppe’s formulation. This is because the kernel of the spatial derivative $\partial = \partial_x$ of the operator product of an arbitrary pair of such Hankel operators can be split into the real matrix product of their respective kernels in the following sense; see Pöppe [55,56]. We include the proof from Doikou et al. [22,21] and Malham [46] for completeness.

**Lemma 1 (Pöppe product)** Assume for each $n \in \mathbb{N}$, $H_n$ are Hankel Hilbert–Schmidt operators with parameter $x$ and $F_n$ are Hilbert–Schmidt operators. Assume further for each $n \in \mathbb{N}$, that the corresponding kernels of $F_n$ are continuous and those of $H_n$ are continuously differentiable. Then, the following Pöppe product rule holds,

$$[F_1 \partial (H_1 H_2) F_2 \partial (H_3 H_4) F_3 \cdots F_{n-1} \partial (H_{n-1} H_n)](y; z; x) = [F_1 H_1(y; 0; x)] [F_2 H_2(0; 0; x)] \cdots [F_{n-1} H_{n-1}(0; 0; x)] [H_n F_n(0; z; x)].$$

**Proof** It is sufficient to establish the result for the case $n = 2$. The general result follows by successive iteration of this case. We use the fundamental theorem of calculus and Hankel properties of $H_1$ and $H_2$. Let $f_1$, $h_1$, $h_2$ and $f_2$ denote the integral kernels of $F_1$, $H_1$, $H_2$ and $F_2$ respectively. By direct computation $[F_1 \partial (H_1 H_2) F_2](y; z; x)$ equals

$$\int_{\mathbb{R}^2} f_1(y, \xi_1; x) \partial_\xi_2 \left( h_1(\xi_1 + \xi_2 + x) h_2(\xi_2 + \xi_3 + x) \right) f_2(\xi_3, z; x) \, d\xi_3 \, d\xi_2 \, d\xi_1$$

$$= \int_{\mathbb{R}^2} f_1(y, \xi_1; x) \partial_\xi_2 \left( h_1(\xi_1 + \xi_2 + x) h_2(\xi_2 + \xi_3 + x) \right) f_2(\xi_3, z; x) \, d\xi_3 \, d\xi_2 \, d\xi_1$$

$$= \int_{\mathbb{R}^2} f_1(y, \xi_1; x) h_1(\xi_1 + x) h_2(\xi_3 + x) f_2(\xi_3, z; x) \, d\xi_3 \, d\xi_1.$$
\[
\begin{align*}
&= \int_{\mathbb{R}^-} f_1(y, \xi_1; x) h_1(\xi_1 + x) \, d\xi_1 \cdot \int_{\mathbb{R}^-} h_2(\xi_3 + x) f_2(\xi_3, z; x) \, d\xi_3 \\
&= ([F_1 H_1](y, 0; x)) ([H_2 F_2](0, z; x)),
\end{align*}
\]
giving the result.

The following lemma and theorem, based on the Leibniz rule, are key to the abstract formalism we introduce in Section 3. Both results are stated in terms of the operator \( P \), though neither result requires \( P \) to be a Hankel operator.

**Lemma 2 (Inverse operator Leibniz rule)** Suppose the operator \( P \) depends on a parameter with respect to which we wish to compute derivatives. Further suppose \( U := (\text{id} - P)^{-1} \) exists. Then we observe

\[
U \equiv \text{id} + UP \equiv \text{id} + PU,
\]
and also that \( \partial U \equiv U(\partial P)U \), and more generally that

\[
\partial^n U \equiv \sum_{k=0}^{n-1} \binom{n}{k} (\partial^k U)(\partial^{n-k} P)U \equiv \sum_{k=0}^{n-1} \binom{n}{k} U(\partial^{n-k} P)(\partial^k U).
\]

**Proof** The first three identities are straightforward. The final identity follows by applying the Leibniz rule to the first pair of identities and using \( U := (\text{id} - P)^{-1} \).

Similar identities to those in Lemma 2 were derived by Pöppe [55,56]. For any \( n \in \mathbb{N} \), let \( C(n) \) denote the set of all compositions of \( n \). The real-valued coefficients \( \chi(a_1 a_2 \cdots a_k) \) given in the following theorem are essentially particular products of Leibniz coefficients, they are given in Definition 5.

**Theorem 1 (Signature operator expansion)** For any non-negative integer \( k \), set \( P_k := \partial^k P \) and \( U_k := \partial^k U \). Then for any \( n \in \mathbb{N} \) we have

\[
U_n = \sum \chi(a_1 a_2 \cdots a_k) UP_{a_1} U_{a_2} U \cdots U_{a_k} U,
\]
where the sum is over all compositions \( a_1 a_2 \cdots a_k \in C(n) \).

We prove this result, which is a direct consequence of Lemma 2, in a more abstract context in the next section; see Theorem 3. Indeed, in Sections 3 and 4 we pull back the definitions and results above to two natural abstract algebras.

Let us now outline the solution procedure based on Pöppe’s Hankel operator approach we employ herein for the non-commutative Korteweg–de Vries hierarchy, and its relation to Grassmannian flows. Solutions to this linear system generate solutions to the target nonlinear partial differential equations from the non-commutative potential Korteweg–de Vries hierarchy or any linear combination of fields from the hierarchy.

**Definition 3 (Linear operator system)** Suppose the linear operators \( P = P(x, t) \) and \( G = G(x, t) \) satisfy the pair of linear equations:

\[
\partial_t P = d(\partial)P \quad \text{and} \quad P = G(\text{id} - P),
\]
where \( d(\partial) = \mu_3 \partial^3 + \mu_5 \partial^5 + \mu_7 \partial^7 + \cdots \) and each \( \mu_n \), with \( n \geq 3 \) odd, is a constant real parameter.
This system of equations for the unknown operator $G$ is linear since $P = P(x, t)$ is the solution of a linear partial differential equation and then $G = G(x, t)$ is the solution of a linear Fredholm equation. Let us now address the existence and uniqueness of such a solution operator in the class of Hilbert–Schmidt operators. Suppose the square matrix-valued function $p = p(x, t)$ satisfies the linear partial differential equation,

$$
\partial_t p = d(\partial)p,
$$

with $p(x, 0) = p_0(x)$ and $p_0$ a given square matrix-valued function. For $w : \mathbb{R} \rightarrow \mathbb{R}_+$, let $L^2_w$ denote the space of real, square matrix-valued functions $f$ on $\mathbb{R}$ whose $L^2$-norm weighted by $w$ is finite, i.e.

$$
\|f\|_{L^2_w} := \int_{\mathbb{R}} \text{tr} \left( f^T(x) f(x) \right) w(x) \, dx < \infty,
$$

where $f^T$ denotes the transpose of $f$ and ‘$\text{tr}$’ is the trace operator. Let $W : \mathbb{R} \rightarrow \mathbb{R}_+$ denote the function $W : x \mapsto 1 + x^2$. Further let $H$ denote the Sobolev space of real, square matrix-valued functions who themselves, as well as derivatives $\partial$ to all orders of them, are square-integrable. Doikou et al. [21] Lemma 3.1 establish if $p_0 \in H \cap L^1_{W}$, then $p \in C^\infty([0, \infty); H \cap L^1_{W})$ and the corresponding Hankel operator $P = P(t)$ whose kernel is $p$, is Hilbert–Schmidt valued. For Hilbert–Schmidt valued operator $P$ we define the regularised Fredholm determinant to be,

$$
det_2(id - P) := \exp \left( -\sum_{k \geq 2} \frac{1}{k} \text{tr} \left( P^k \right) \right).
$$

See Simon [64]. The operator $id - P$ is invertible if and only if $\det_2(id - P) \neq 0$. Now assume as just above, $p_0 \in H \cap L^1_{W}$, and further that the Hilbert–Schmidt Hankel operator $P_0$ corresponding to the kernel $p_0$ is such that $\det_2(id - P_0) \neq 0$. Then from Doikou et al. [21] Lemma 3.2 we know, there exists a time $T > 0$, such that for all $t \in [0, T]$ and $x \in \mathbb{R}$, we know $\det_2(id - P(x, t)) \neq 0$, where $P = P(x, t)$ is the Hilbert–Schmidt Hankel operator whose kernel is $p = p(y + x; t)$ given just above, and $P(x, 0) = P_0(x)$. Further, there exists a unique square matrix-valued $g \in C^\infty([0, T]; C^\infty(\mathbb{R}^2 \times \mathbb{R}_+))$ which satisfies the linear Fredholm equation,

$$
p(y + z + x; t) = g(y, z; x, t) - \int_{-\infty}^{0} g(y, \xi; x, t) p(\xi + z + x, t) \, d\xi,
$$

corresponding to the operator equation $P = G(id - P)$ for $G = G(x, t)$.

Let us now briefly outline why the solution flow for $G$ given in Definition 3 is a Fredholm Grassmannian flow; many more details can be found in Beck et al. [10;11], Doikou et al. [21] Sec. 2.3 and Doikou et al. [22]. Detailed introductions and examples of Fredholm Grassmannian manifolds are given in Pressley and Segal [59], Abbondandolo and Majer [1] and Andrucho and Larotonda [7]. Suppose $H$ is a given separable Hilbert space. The Fredholm Grassmannian of all subspaces of $H$ that are comparable in size to a given closed subspace $V \subset H$ is defined as follows; see Segal and Wilson [63].

**Definition 4 (Fredholm Grassmannian)** Let $H$ be a separable Hilbert space with a given decomposition $H = V \oplus V^\perp$, where $V$ and $V^\perp$ are infinite dimensional closed subspaces. The Grassmannian $\text{Gr}(H, V)$ is the set of all subspaces $W$ of $H$ such that:
(i) The orthogonal projection $pr : W \to V$ is a Fredholm operator, indeed it is a Hilbert–Schmidt perturbation of the identity; and

(ii) The orthogonal projection $pr : W \to V^\perp$ is a Hilbert–Schmidt operator.

We observe that as $H$ is separable, any element in $H$ has a representation on a countable basis. Thus, for example, we can represent that element via the sequence of coefficients which are associated with each of the basis elements. Given an independent set of such sequences in $H = V \oplus V^\perp$ which span $V$, suppose we record them as columns in the infinite matrix

$$W = \begin{pmatrix} \id + Q \\ P \end{pmatrix}.$$ 

In this representation, we suppose each column of $\id + Q$ lies in $V$ and each column of $P$ lies in $V^\perp$. Further, assume when we constructed $\id + Q$, we ensured it was a Fredholm operator on $V$ with $Q \in J_2(V;V)$, where $J_2(V;V)$ is the class of Hilbert–Schmidt operators from $V$ to $V$, equipped with the norm $\|Q\|_2^2 = \tr Q^\dagger Q$ where ‘tr’ is the trace operator. See Simon [64] for more details. In addition, assume we constructed $P$ to ensure $P \in J_2(V;V^\perp)$, the space of Hilbert–Schmidt operators from $V$ to $V^\perp$.

Let $\mathbb{W}$ denote the subspace of $H$ represented by the span of the columns of $W$. Let $V_0 \cong V$ denote the canonical subspace of $H$ with the representation

$$V_0 = \begin{pmatrix} \id \\ O \end{pmatrix},$$ 

where $O$ is the infinite matrix of zeros. Assume that $\det_2(\id + Q) \neq 0$. Then the projections $pr : \mathbb{W} \to V_0$ and $pr : \mathbb{W} \to V_0^\perp$ respectively generate

$$W^\parallel = \begin{pmatrix} \id + Q \\ O \end{pmatrix} \quad \text{and} \quad W^\perp = \begin{pmatrix} O \\ P \end{pmatrix}.$$ 

The subspace of $H$ represented by the span of the columns of $W^\parallel$ naturally coincides with the subspace $V_0$. Indeed, the transformation $(\id + Q)^{-1} \in \GL(V)$ transforms $W^\parallel$ to $V_0$. Under this transformation the representation $W$ for $\mathbb{W}$ becomes

$$\begin{pmatrix} \id \\ G \end{pmatrix},$$ 

where $G = P(\id + Q)^{-1}$. In fact, any subspace $\mathbb{W}$ that can be projected onto $V_0$ can be represented in this way and vice-versa. In this representation the operators $G \in J_2(V;V^\perp)$ parametrise all the subspaces $\mathbb{W}$ that can be projected on to $V_0$. If $\det_2(\id + Q) = 0$, then the projection above is not possible. In this instance we simply choose a different representative coordinate chart/patch. It is always possible to choose a subspace $V_0' \cong V$ of $H$ such that the projection $\mathbb{W} \to V_0'$ is an isomorphism; see Pressley and Segal [59, Prop. 7.1.6]. For further details on coordinate patches see Beck et al. [10] and Doikou et al. [22] and Doikou et al. [21]. In conclusion we observe that the flow for $G$ given in Definition [3] taking $Q$ to be $-P'$ and under the assumptions stated for $p_0$ directly after, including the regularised determinant restriction, represents a Fredholm Grassmannian flow.
3 Hilbert–Schmidt operator and concatenation algebras

Let $C := \cup_{n \geq 0} C(n)$ denote the set of all compositions. Herein we consider an abstract version of the operator algebra generated by a parameter-dependent operator $P$, its derivatives with respect to the parameter, and by $U := (id - P)^{-1}$ as well as the derivatives of $U$ with respect to the parameter. With $P_n := \partial^n P$ and $U_n := \partial^n U$ for all non-negative integers $n$, let $\mathcal{P}$ denote the alphabet $\{P_n\}_{n \geq 0}$ and $\mathcal{U}$ denote the alphabet $\{U_n\}_{n \geq 0}$. Then we denote by $\mathcal{R}(\mathcal{P} \cup \mathcal{U})$, the Hilbert–Schmidt operator algebra over $\mathbb{R}$, generated by $\mathcal{P}$ and $\mathcal{U}$. This algebra is unital with the identity operator as the unit element. One natural approach to generating an abstract version of the algebra $\mathcal{R}(\mathcal{P} \cup \mathcal{U})$ is as follows.

Let $\mathbb{N}$ denote the set of non-negative integers $n$. We distinguish such a set of integers for the following reason. In the first stage of the abstraction we essentially replace the operators $P_n$ by integers $n$, and the operators $U_n$ by the integers $n$, and construct a concatenation algebra of words generated by the alphabets $\mathbb{N}$ and $\mathbb{N}$. We denote by $\mathcal{R}(\mathbb{N} \cup \mathbb{N})$ the non-commutative concatenation polynomial algebra over $\mathbb{R}$ generated words constructed from the alphabet $\mathbb{N} \cup \mathbb{N}$; see Reutenauer [60]. The neutral element is the empty word $\emptyset$, so that the concatenation of any word from $\mathcal{R}(\mathbb{N} \cup \mathbb{N})$ with $\emptyset$ and vice-versa, just generates the original word. Hence $\mathcal{R}(\mathbb{N} \cup \mathbb{N})$ is a unital algebra. The neutral element $\emptyset$ plays the role equivalent to that of the identity operator in $\mathcal{R}(\mathcal{P} \cup \mathcal{U})$. Indeed, from their definitions we observe that $\mathcal{R}(\mathbb{N} \cup \mathbb{N})$ and $\mathcal{R}(\mathcal{P} \cup \mathcal{U})$ are isomorphic. We can also define on $\mathcal{R}(\mathbb{N} \cup \mathbb{N})$ a derivation operation $\delta$ that represents the abstraction of the derivative operator $\partial$ acting on the operators $P$ and $U$. Indeed operator $\delta$: $\mathcal{R}(\mathbb{N} \cup \mathbb{N}) \to \mathcal{R}(\mathbb{N} \cup \mathbb{N})$ is the linear operator such that for any letter $n \in \mathbb{N}$ in $\mathcal{R}(\mathbb{N} \cup \mathbb{N})$, $\delta: n \mapsto (n + 1)$, and for any letter $n \in \mathbb{N}$ in $\mathcal{R}(\mathbb{N} \cup \mathbb{N})$, $\delta: n \mapsto (n + 1)$. Furthermore for any word $a_1a_2a_3\cdots a_n \in \mathcal{R}(\mathbb{N} \cup \mathbb{N})$, where the letters $a_1, a_2, \ldots, a_n$ can be either from $\mathbb{N}$ or $\mathbb{N}$, the operator $\delta$ satisfies the Leibniz derivation property:

$$\delta(a_1a_2a_3\cdots a_n) = \sum_{k=1}^{n} a_1a_2\cdots a_{k-1}(\delta a_k)a_{k+1}\cdots a_n.$$ 

The first crucial result of this section is that each of the letters in $\mathbb{N}$ has a linear expansion in terms of monomials of the form $a_1a_20\cdots a_n0$ with $a_1, a_2, \ldots, a_n \in \mathbb{N}$ only. Note here the letter $0$ corresponds to $U_0 \equiv U$. This is the abstract version of the statement of Theorem [4]. We state and prove this result after the following definition. We denote by $\mathbb{N}^*$, the free monoid of words on $\mathbb{N}$.

**Definition 5 (Signature character)** Let $a_1a_2\cdots a_n$ be a word from $\mathbb{N}^*$. We associate with any such word the signature character $\chi: \mathbb{N}^* \to \mathbb{R}$, given by

$$\chi(a_1a_2\cdots a_n) := \left(a_1 + \cdots + a_n \atop a_n\right) \left(a_1 + a_2 + \cdots + a_n \atop a_1 \right) \cdots \left(a_{n-1} + a_k \atop a_{n-1}\right) \left(a_n \atop a_n\right),$$

where each of the factors shown on the right is a Leibniz coefficient, so that for example, the penultimate factor is $a_{n-1} + a_n$ choose $a_{n-1}$.

Let us introduce some notation. Given a word $w = a_1a_2\cdots a_n$ generated using letters $a_1, a_2, \ldots, a_n$ from $\mathbb{N}$, let $w_0$ denote the corresponding word in $\mathcal{R}(\mathbb{N} \cup \mathbb{N})$ of the form

$$w_0 := a_1a_20\cdots 0a_n0,$$

i.e. a single letter $0$ exists between each letter $a_1, a_2, \ldots, a_n$ from $\mathbb{N}$, and at each end.
Theorem 2 (Signature operator expansion) Given any \( n \in \mathbb{N} \), we have, in the algebra \( \mathbb{R}(\mathbb{N} \cup \mathbb{N}) \), for any corresponding integer letter \( n \) from \( \mathbb{N} \) that:

\[
\mathbf{n} = \sum_{w \in C(n)} \chi(w) \cdot w_0.
\]

Proof We begin by recreating some of the basic results for the operators \( P \) and \( U \) from \( \mathbb{R}(\mathbb{P} \cup \mathbb{U}) \) in Lemma 2 in terms of the corresponding elements from \( \mathbb{R}(\mathbb{N} \cup \mathbb{N}) \). With the letter 0 corresponding to \( P \), the letter 0 corresponding to \( U \) and with \( \emptyset \) representing the neutral element corresponding to the identity \( id \) in \( \mathbb{R}(\mathbb{P} \cup \mathbb{U}) \), we have

\[
0 \equiv (\emptyset - 0)^{-1} \quad \text{and thus} \quad 0 \equiv \emptyset + 00 \equiv \emptyset + 00.
\]

These correspond to the first set of statements in Lemma 2. Next using that \( \delta(n) = (n + 1) \), we apply \( \delta \) several times to the identity \( \emptyset \equiv \emptyset + 00 \). Thus we have

\[
1 \equiv \begin{pmatrix} 1 \end{pmatrix} \cdot 10 + \begin{pmatrix} 1 \end{pmatrix} \cdot 01 \quad \leftrightarrow \quad 1 \equiv \begin{pmatrix} 1 \end{pmatrix} \cdot 010.
\]

This establishes the statement of the lemma for 1. One further iteration of this procedure is insightful. Applying \( \delta^2 \) to the identity \( \emptyset \equiv \emptyset + 00 \), we observe

\[
2 \equiv \begin{pmatrix} 2 \end{pmatrix} \cdot 20 + \begin{pmatrix} 2 \end{pmatrix} \cdot 11 + \begin{pmatrix} 2 \end{pmatrix} \cdot 02 \quad \leftrightarrow \quad 2 \equiv \begin{pmatrix} 2 \end{pmatrix} \cdot 020 + \begin{pmatrix} 2 \end{pmatrix} \cdot 011
\]

Substituting for 1 into the right-hand side, we find

\[
2 \equiv \begin{pmatrix} 2 \end{pmatrix} \cdot 020 + \begin{pmatrix} 2 \end{pmatrix} \cdot \begin{pmatrix} 1 \end{pmatrix} \cdot 010.
\]

This corresponds to the case \( n = 2 \) in the statement of the lemma. We now focus on the overall proof which proceeds by induction. We have seen the statement of the lemma holds for \( n = 1 \). Assume the statement holds for all \( 1 \leq k \leq n \). We now show it holds for \( (n + 1) \), i.e. that we have

\[
(n + 1) = \sum_{w \in C(n+1)} \chi(w) \cdot w_0.
\]

We apply \( \delta^{n+1} \) to the identity \( \emptyset \equiv \emptyset + 00 \). Using the Leibniz rule and combining the final term \( 0(n+1)0 \) on the right with the term \( \delta^{n+1}0 \equiv (n+1) \) on the left, in the same manner as we did above, we obtain

\[
(n + 1) \equiv \begin{pmatrix} n+1 \end{pmatrix} \cdot 0(n+1)0 + \begin{pmatrix} n+1 \end{pmatrix} \cdot 0n1 + \begin{pmatrix} n+1 \end{pmatrix} \cdot 0(n-1)2
\]

\[
+ \begin{pmatrix} n+1 \end{pmatrix} \cdot 0(n-2)3 + \cdots + \begin{pmatrix} n+1 \end{pmatrix} \cdot 02(n-1) + \begin{pmatrix} n+1 \end{pmatrix} \cdot 01n.
\]

We have assumed that for each \( k \) such that \( 2 \leq k \leq n \), we have \( k \) is the appropriate linear combination of all the compositions of \( k \). The first term on the right above involving \( 0(n+1)0 \) is the only composition of \( n+1 \) starting with the digit \( n+1 \) — ignoring the 0’s. This term on the right above matches the term in \( (n + 1) \) corresponding to the single digit composition of \( n+1 \), with the correct coefficient \( \chi(n+1) \). Now consider the next term on the right involving \( 0n1 \). If we substitute for 1 we obtain the the only
composition of \( n + 1 \) starting with the digit \( n \)—again ignoring the \( 0 \)'s—namely the term \( 0a010 \). The coefficient is

\[
\binom{n + 1}{n} \binom{1}{1} = \chi(n1).
\]

Indeed, for each \( 1 \leq k \leq n \), consider the term involving \( 0(n + 1 - k)k \) on the right in the Leibniz expansion above. We substitute for \( k \) which is the appropriate linear combination of all compositions of \( k \). We observe this exhausts all the possible compositions of \( n + 1 \) that start with the digit \( n + 1 - k \)—ignoring the \( 0 \)'s. Now consider the corresponding coefficients. Consider a generic term \( w_0 \) in \( k \), for which the \( 0 \)-stripped version of \( w_0 \), namely \( w \), is a composition of \( k \). The coefficient of \( w_0 \) is \( \chi(w) \). Its coefficient in the Leibniz expansion above would be

\[
\binom{n + 1}{n + 1 - k} \chi(w) = \chi((n + 1 - k)w).
\]

This last result follows from Definition 5 for \( \chi(w) \): for any word \( w \) which is a composition of \( k \) and for which for a given letter \( a \), the word \( aw \) is a composition of \( n + 1 \), then \( a = n + 1 - k \) and the result above holds. This completes the proof. \( \square \)

The significance of Theorem 2 is that every element \( n \in \mathbb{R}(\mathbb{N} \cup \mathbb{N}) \) for \( n \in \mathbb{N} \) can be expressed as a linear combination of monomials of the form \( 0a_10a_20 \cdots 0a_n0 \) with all the \( a_k \in \mathbb{N} \), i.e. from the integer alphabet only, not involving \( \mathbb{N} \). This means that we can use monomials of this form as a basis for \( n \in \mathbb{R}(\mathbb{N} \cup \mathbb{N}) \). Indeed it thus is sufficient for us to consider the concatenation algebra \( n \in \mathbb{R}(\mathcal{C}_0) \) of monomials of the form \( 0a_10a_20 \cdots 0a_n0 \) from \( \mathcal{C}_0 \), the set of all compositions with the letter 0 squeezed between each digit as well as being added to each end.

4 Hankel kernel and Pöppe algebras

We have seen that, as a result of Theorem 3, it is sufficient for us to consider the concatenation algebra \( n \in \mathbb{R}(\mathcal{C}_0) \) of monomials of the form \( 0a_10a_20 \cdots 0a_n0 \) where \( a_1a_2 \cdots a_n \in \mathcal{C} \). This means instead of the Hilbert–Schmidt operator algebra \( \mathbb{R}(\mathcal{P} \cup \mathcal{U}) \), it is sufficient for us to consider the Hilbert–Schmidt operator algebra \( \mathbb{R}(\mathcal{U}) \) generated by operator monomials of the form \( UP_{a_1}UP_{a_2}UP_{a_3} \cdots UP_{a_n}UP \). Indeed we could have deduced this last statement directly from Theorem 4. Our goal herein is to consider the algebra of Hilbert–Schmidt Hankel kernels generated by applying the the kernel bracket \([\cdot \cdot\cdot]\) to such operator monomials. To this end, let \( \mathbb{R}[\mathcal{U}] \) denote the non-commutative polynomial matrix algebra generated by \textit{kernel monomials} of the form \( [UP_{a_1}UP_{a_2}UP \cdots UPP_{a_n}UP] \) where \( a_1a_2 \cdots a_n \in \mathcal{C} \). We call this the \textit{Hankel kernel algebra}. We endow this algebra with a product equivalent to the Pöppe product in Lemma 4 for such kernel monomials as follows.

Lemma 3 (Pöppe product for kernel monomials) \( \text{Suppose } u = a_1 \cdots a_k \in \mathcal{C} \) and \( v = b_1 \cdots b_l \in \mathcal{C} \) while \( a, b \in \mathbb{N} \). Let \( m_u \) and \( m_v \) denote the respective kernel monomials \( m_u := UP_{a_1}UP_{a_2}UP \cdots UP_{a_k}UP \) and \( m_v := UP_{b_1}UP_{b_2}UP \cdots UP_{b_l}UP \). Then we have

\[
[m_uP_0U][UP_bP_0U] = [m_uP_{a+1}UP_bP_0U] + [m_uP_{b+1}UP_bP_0U] + 2[m_uP_0UP_bUP_mP_0].
\]
Proof Using that \( \partial U \equiv \text{id} + UP \equiv \text{id} + PU \), the P"oppe product rule from Lemma 1 and that \( \partial U = (\partial PU)U \), we observe

\[
[m_a P_a U][UP_b m_v] = [m_a P_a + m_a P_a UP][PU_P b_m_v + P_b m_v] \\
= [m_a P_a][P_a m_v] + [m_a P_a][PU_P b_m_v] \\
+ [m_a P_a U][P_b m_v] + [m_a P_a U][PU_P b m_v] \\
= [m_a \partial (P a P)U P b_m_v] + [m_a \partial (P a P)U P b_m_v] \\
+ [m_a P a U \partial (P P b) m_v] + [m_a P a U \partial (P P b)U P b m_v] \\
= [m_a \partial (P a P)U P b m_v] + [m_a P a U \partial (P P b) m_v] \\
+ [m_a \partial (P a P)U P b m_v] + [m_a P a U \partial (P P b)U P b m_v] \\
= [m_a \partial (P a P)U P b m_v] + [m_a P a U \partial (P P b) m_v] \\
+ [m_a \partial (P a P)U P b m_v] + [m_a P a U \partial (P P b)U P b m_v] \\
= [m_a \partial (P a P)U P b m_v] + [m_a \partial (P a P)U P b m_v] \\
+ 2 [m_a P a U \partial (P P b)U P b m_v],
\]

which establishes the result. \( \square \)

Remark 1 Recall from Lemma 1, we implicitly interpret multiple kernel products of the form \([-\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\cdot\……
the product \((aa)∗(bv)\) can be interpreted analogously. In this \(0\)-padded abstraction, we interpret \(0\) as corresponding to the pseudo-passive operator \(‘U’\) in the kernel monomials \([UP_a UP_{a_2} U \cdots UP_{a_n} U]\). With this interpretation there is a one-to-one correspondence between the monomials \(0a_1 0a_2 \cdots 0a_n 0\) and \([UP_a UP_{a_2} U \cdots UP_{a_n} U]\). The Pöpppe products given in Lemma \(3\) and Definition \(7\) are mirror images of each other and the algebras generated by the corresponding sets of monomials are seen to be isomorphic. However in practice, keeping the pseudo-passive \(0\)’s is cumbersome, and not required.

Herein we use the stripped abstraction of \(\mathbb{R}([P_U])\), replacing it by its isomorphic equivalent \(\mathbb{R}(C)_∗\). The coding we use is to replace the kernel monomials \([UP_a UP_{a_2} U \cdots UP_{a_n} U]\) by the \(0\)-stripped compositions \(a_1 \cdots a_n\) and use the Pöpppe product for compositions in Definition \(7\) in place of the Pöpppe product in Lemma \(3\). Hence, to emphasise, \(\mathbb{R}(C)_∗\) is the non-commutative polynomial algebra over \(\mathbb{R}\) generated by the set of compositions \(C\) with the product being the Pöpppe product in Definition \(7\). Further, from their definitions, we observe that \(\mathbb{R}(C)_∗\) and \(\mathbb{R}([P_U])\) are isomorphic.

**Remark 3** Note we do not require \(\mathbb{R}([P_U])\) and \(\mathbb{R}(C)_∗\) to be unital algebras here nor in the proof of our main result Theorem \(3\) in Section \(3\). However when we construct the signature co-algebra in Section \(3\) it is useful to introduce the empty composition \(ν ∈ \mathbb{R}(C)_∗\) with the property \(ν ∗ w = w ∗ ν = w\) for any composition \(w\).

The abstract versions of the signature expansions \([Un]\) for kernels in \(\mathbb{R}([P_U])\) in Definition \(6\) are the following signature expansions in \(\mathbb{R}(C)_∗\).

**Definition 8 (Signature expansions)** For any \(n ∈ \mathbb{N}\) we define the following linear expansions \(ν ∈ \mathbb{R}(C)_∗\): \(n := \sum_{w ∈ C(n)} χ(w) · w\).

**Remark 4 (Notation)** Note, in Definition \(8\) and hereafter, we use \(1, 2, \ldots, n\) and so forth to denote linear signature expansions in \(\mathbb{R}(C)_∗\) associated with the kernel algebra \(\mathbb{R}([P_U])\), i.e. as linear expansions in terms of composition monomials in \(\mathbb{R}(C)_∗\). Hitherto, see for example Theorem \(2\) we used \(n\) to denote signature expansions in \(\mathbb{R}(N ∪ N)\) associated with the operator algebra \(\mathbb{R}([P ∪ U])\).

**Remark 5 (Convention)** In Definition \(8\) we observe a notation convention we have used thus far for expansions in \(\mathbb{R}(C)_∗\) and \(\mathbb{R}(C)_0\) and which we use hereafter. That is, we distinguish between the real coefficients associated with combinatorial elements such as \(a_1 \cdots a_k \in C\) by separating them by a ‘,’ with the real-valued coefficients residing on the left and the compositions residing on the right. So for example if we write \(3 · 221\), this is to be interpreted as \(3\) copies of the composition element ‘221’.

As for the abstract operator algebra in Section \(3\) we can define on \(\mathbb{R}(C)_∗\) a derivation operation \(\partial\) as follows.

**Definition 9 (Derivation in composition kernel algebra)** Given any monomial \(w ∈ \mathbb{R}(C)_∗\) with \(w = a_1 a_2 \cdots a_n\) we define the derivation \(\partial\) of the composition \(w\) to be the linear expansion: \(\partial w = \sum_{k=1}^{n} a_1 a_2 \cdots a_{k-1} (a_{k+1} + 1) a_{k+1} \cdots a_n + \sum_{k=1}^{n+1} a_1 a_2 \cdots a_{k-1} 1 a_k \cdots a_n\).

In the second sum, the \(k = 1\) and \(k = n + 1\) cases correspond to a ‘1’ being appended, respectively, to the front and then end of the composition \(a_1 a_2 \cdots a_n\).
Example 1 For any two arbitrary single letter compositions $a, b \in \mathbb{R} \langle C \rangle_*$, we have $\partial(ab) = 1ab + (a + 1)b + ab + (b + 1) + ab1$. More generally, for any two arbitrary compositions $ua, bv \in \mathbb{R} \langle C \rangle_*$, where we distinguish the final and beginning single letters $a$ and $b$ as shown, we have $\partial(ua) = (\partial u)abv + u(\partial(ab))v + uab(\partial v)$.

Remark 6 This can be viewed as the result of the derivation of the 0-padded version $w_0$ of $w$ using the Leibniz rule. In other words the definition above is just the abstract encoding of applying the derivation operation $\partial$ to $[UPa_1 UPa_2 U \cdots UPa_n U]$. Note the linear operations $\partial$ and $[\cdot]$ commute.

Just as for the abstract operator algebra, for the abstract kernel algebra here we have the following property.

Lemma 4 For any signature expansion $n \in \mathbb{R} \langle C \rangle_*$, the action of the derivation operation is $\partial: n \mapsto (n + 1)$, i.e. we have $\partial(n) = (n + 1)$.

Proof The derivation $\partial(n)$ corresponds to the signature expansion in Definition 8 with all the compositions $w \in C(n)$ replaced by $\partial(w)$. Using that $\mathbb{R} \langle C \rangle_*$ and $\mathbb{R} \langle \mathbb{P}_U \rangle$ are isomorphic, this corresponds to the signature expansion for kernels in Definition 8 with each term $[UPa_1 UPa_2 U \cdots UPa_n U]$ therein replaced by $\partial[UPa_1 UPa_2 U \cdots UPa_n U]$. We know this expansion, by applying the differential operator $\partial$ and the kernel bracket operator $[\cdot]$ to the result of Theorem 1, corresponds to $\partial[U_n]$ which equals $[U_{n+1}]$. The signature kernel expansion for $[U_{n+1}]$ corresponds to that for $(n + 1)$, by comparing term for term, their signature expansions in terms of compositions $w \in C(n + 1)$. □

5 Signature co-algebra

Herein we motivate and introduce the signature co-algebra $S$. Though not strictly necessary to prove our main result, our impetus for constructing the co-algebra will be useful in helping to elucidate our main result in Section 6. To motivate our construction of the signature co-algebra consider the product of two arbitrary signature expansions $m$ and $n$ in $\mathbb{R} \langle C \rangle_*$:

$$m \star n = \sum \chi(u)\chi(v) \cdot u \star v,$$

where the sum is over all $u \in C(m)$ and $v \in C(n)$, and the product is naturally the Pöppe product for compositions. For the real-valued coefficients $\chi(u)$ and $\chi(v)$ we have $\chi(u)\chi(v) \equiv \chi(v)\chi(u)$ for any pair of compositions $u, v \in C$; recall the signature character map $\chi$ from Definition 2. In the product above, we can expand the products $u \star v$ into linear combinations of compositions using Definition 7 for the Pöppe product. However we observe that, in the sum on the right above for the product $m \star n$, once we have expanded all the terms $u \star v$ for all $u \in C(m)$ and $v \in C(n)$, it might be useful to explicitly preserve the generators of specific compositions that appear in the sum on the right; and we should extend this idea to higher degree versions of such products.

Example 2 Consider the linear combination $c_{21} \cdot 2 * 1 + c_{12} \cdot 1 * 2$ for some real constants $c_{21}$ and $c_{12}$. The product $2 * 1$ equals

$$2 * 1 = (\chi(2) \cdot 2 + \chi(11) \cdot 11) * (\chi(1) \cdot 1) = \chi(2)\chi(1) \cdot 2 * 1 + \chi(11)\chi(1) \cdot 11 * 1 = \chi(2)\chi(1) \cdot (31 + 22 + 2 \cdot 211) + \chi(11)\chi(1) \cdot (121 + 112 + 2 \cdot 1111)$$
The reason for introducing the special element ‘q’ is now explained. When we encode the coefficients in this tensorially manner we need a mechanism to record those compositions that are generated by the third term in the Pöppe product of two appropriate compositions. We use this encoding explicitly when proving our main theorem.

\[
\chi(2)\chi(1) \cdot (31 + 22) + \chi(2)\chi(1)\chi(q) \cdot (211) + \chi(11)\chi(1) \cdot (121 + 112) + \chi(11)\chi(1)\chi(q) \cdot (1111),
\]

and so we can deduce the product \(1 \ast 2\) equals

\[
1 \ast 2 = \chi(1)\chi(2) \cdot (13 + 22) + \chi(1)\chi(2)\chi(q) \cdot (112) + \chi(1)\chi(11) \cdot (121 + 211) + \chi(1)\chi(11)\chi(q) \cdot (1111).
\]

Note we distinguish a special element ‘q’ to indicate it results from the third term in the Pöppe product on compositions, so \(\chi(q) = 2\). For this example, think of the symbol ‘(q)’ simply as a proxy for the real coefficient ‘2’. So if we expand the linear combination \(c_{21} \ast 1 + c_{12} \ast 1 \ast 2\), by preserving the signature character coefficients in the manner above, we know the composition 31 was the result of the product 2 \ast 1, because the coefficient of 31 is \(\chi(2)\chi(1)\). On the other hand, the composition 211 which would appear with the coefficient \(c_{21}\chi(2)\chi(1)\chi(q) + c_{12}\chi(1)\chi(11)\) was generated as follows.

It was generated by the third term in the product 2 \ast 1 coming from the product \(2 \ast 1\), i.e. where a ‘1’ is inserted between the 2 and the 1. That it was this operation that generated this term is indicated by the product of factors \(\chi(2)\chi(1)(q)\), i.e. in particular that it involves the factor \(\chi(q)\). It was also generated in the product 1 \ast (11) coming from the product \(1 \ast 2\), as indicated by the product of factors \(\chi(1)\chi(11)\). In other words, each composition can be associated with a finite product of ‘signatures’.

This example illustrates the usefulness of carefully retaining the coefficients when expanding linear combinations of monomials consisting of Pöppe products of signature expansions. More generally we require the multi-factor version of the product of signature expansions given above of the form:

\[
n_1 \ast n_2 \ast \cdots \ast n_k = \sum \chi(u_1)\chi(u_2)\cdots\chi(u_n) \cdot u_1 \ast u_2 \ast \cdots \ast u_k,
\]

where the sum is over all \(u_1 \in \mathcal{C}(n_1), \ldots, u_k \in \mathcal{C}(n_k)\). Naturally we observe that any composition \(w \in \mathcal{C}\) can result from one or more Pöppe products of the form \(u_1 \ast u_2 \ast \cdots \ast u_k\). It is now apparent that it is convenient to encode \(\chi(u_1)\chi(u_2)\cdots\chi(u_k)\) as \(\chi(u_1 \otimes u_2 \otimes \cdots \otimes u_k)\). Indeed we can assume \(\chi\) to act homomorphically on any such tensor product of compositions so that indeed \(\chi(u_1 \otimes u_2 \otimes \cdots \otimes u_k) \equiv \chi(u_1)\chi(u_2)\cdots\chi(u_k)\).

Thus, for example, we now write:

\[
2 \ast 1 = \chi(2 \otimes 1) \cdot (31 + 22) + \chi(2 \otimes q \otimes 1) \cdot (211) + \chi(11 \otimes 1) \cdot (121 + 112) + \chi(11 \otimes q \otimes 1) \cdot (1111).
\]

The reason for introducing the special element ‘q’ is now explained. When we encode the coefficients in this tensorially manner we need a mechanism to record those compositions that are generated by the third term in the Pöppe product of two appropriate compositions. We use this encoding explicitly when proving our main theorem results in Section 6 always keeping in mind though that expressions of the form \(\chi(u_1 \otimes u_2 \otimes \cdots \otimes u_k)\), where one of the tensorially elements may be ‘q’, are really to interpreted as \(\chi(u_1)\chi(u_2)\cdots\chi(u_k)\) and \(\chi(q)\) just acts as a proxy for the coefficient ‘2’.

The other important idea we extract from our computations above is that it might be useful, for any given composition \(w \in \mathcal{C}\), to determine all the possible pairs of compositions that might have generated it. For example, we observe in Example 2 above that the composition 22 can be generate through 2 \ast 1 and also 1 \ast 2. On the
other hand, the composition 112 can be generated via 11∗1, or via 1∗2 through the third term in the Pöppe product. This can be used as a simple check that we have collated terms correctly, see for example the rows in Table 1 in Section 6. We explore this idea more formally just below. However, the notions we have discussed thusfar are all that we require from this section for our proofs in Sections 6 and 7. At this point, the reader only focused on our main results can safely skip to Section 6.

Remark 7 (Alternative encodings) The encoding above involving the special element ‘q’ is one functional encoding for our main results in Section 6. However, it does not preserve the number of tensorial components. It is perhaps more natural for the number of tensorial components to match the number of Pöppe products associated with generating that particular composition. See Definition 10 and in particular Lemma 5 for de-Pöppe co-product just below to see how this can be achieved.

The rest of this section is devoted to establishing a new co-algebra we call the signature co-algebra. The motivating idea, from the notions just discussed, is to formalise the process of determining all the possible Pöppe products of compositions that generate a given composition. We begin by defining a co-product on C, the de-Pöppe co-product.

Definition 10 (De-Pöppe co-product) For any composition w ∈ C, we define the de-Pöppe co-product ∆(w) of w to be

\[ \Delta(w) = \sum_{u \otimes v \in \mathbb{C}^2} \langle u \ast v, w \rangle \mathbb{C} u \otimes v, \]

where \( \langle \cdot, \cdot \rangle \mathbb{C} \) is the inner product on C defined for any u, v ∈ C by

\[ \langle u, v \rangle \mathbb{C} = \begin{cases} 1 & \text{if } u = v, \\ 0 & \text{if } u \neq v. \end{cases} \]

We give the definition of the signature co-algebra first, and then prove that the signature co-algebra is indeed a co-algebra second.

Definition 11 (Signature co-algebra) We define the signature co-algebra S as the co-algebra \( \mathbb{R} \langle \mathbb{C}^n \rangle \) over \( \mathbb{R} \) constructed from all possible monomials \( s_1 \otimes s_2 \otimes \cdots \otimes s_n \) chosen from \( \mathbb{C}^n \) for all \( n \in \mathbb{N} \). Here \( \mathbb{C}^n \) denotes \( \bigcup_{n \geq 0} \mathbb{C}^n \). We define the co-product \( \Delta : S \rightarrow S \otimes S \) on S to be the de-Pöppe co-product \( \Delta \) in Definition 10. We see from Lemma 5 just below, \( \Delta \) tensorally decomposes any composition w into the sum of all possible composition pairs that produce w via the Pöppe product ‘∗’, including the empty composition \( \nu \); see Remark 5. The co-unit on S which we denote \( \varepsilon : S \rightarrow \mathbb{R} \) is given for any \( w \in \mathbb{C} \) by

\[ \varepsilon(w) = \begin{cases} 1 & \text{if } w = \nu, \\ 0 & \text{if } w \neq \nu. \end{cases} \]

We can derive an explicit formula for the de-Pöppe product as follows. To achieve this, setting \( 0 = 0^{-1}(1) \), the following map \( \theta : \mathbb{C} \cup \{0\} \rightarrow \mathbb{C} \) proves useful. For any composition \( w \in \mathbb{C} \), we define:

\[ \theta : w \mapsto w; \quad \theta : 0w \mapsto w; \quad \theta : w0 \mapsto w; \quad \theta : \nu \mapsto \nu \quad \text{and} \quad \theta : 0 \mapsto 0 \cdot \nu. \]

The coefficient in the image in the final case is 0 ∈ \( \mathbb{R} \), so ultimately the term is zero.
Lemma 5 (De-Poppe co-product formula) The co-product $\Delta$ can be characterised as follows. For any composition $a_1 a_2 \cdots a_n$ we have $\Delta(a_1 a_2 \cdots a_n)$ is given by:

$$(\theta \otimes \theta) \circ \left( \sum_{k=0}^{n} (a_1 \cdots a_{k-1}(a_k-1) \otimes a_{k+1} \cdots a_n + a_1 \cdots a_k \otimes (a_{k+1}-1)a_{k+2} \cdots a_n) \right),$$

where the $k = 0$ and $k = n$ terms are $\nu \otimes a_1 a_2 \cdots a_n$ and $a_1 a_2 \cdots a_n \otimes \nu$, respectively.

Proof For any composition $a_1 a_2 \cdots a_n \in \mathcal{C}$, consider the argument of the map $\theta \otimes \theta$ given in the statement of the theorem. This is equivalent to the re-written form:

$$\nu \otimes a_1 a_2 \cdots a_n + (a_1 - 1) \otimes a_2 \cdots a_n + a_1 \otimes (a_2 - 1) a_3 \cdots a_n + a_1(a_2 - 1) \otimes a_3 \cdots a_n + a_1 a_2 \otimes (a_3 - 1) a_4 \cdots a_n + \cdots + a_1 \cdots (a_{n-1} - 1) \otimes a_n + a_1 \cdots a_{n-1} \otimes (a_n - 1) + a_1 a_2 \cdots a_n \otimes \nu.$$

The first and last terms account for the possibility $w = a_1 a_2 \cdots a_n$ could be generated by $\nu \ast w$ and $w \ast \nu$, respectively. In general we observe that for $k = 1, \ldots, n-1$ the term $a_1 \cdots a_k a_{k+1} \cdots a_n$ could be generated by either $a_1 \cdots a_{k-1}(a_k - 1) \ast a_{k+1} \cdots a_n$ or $a_1 \cdots a_k \ast (a_{k+1} - 1)a_{k+2} \cdots a_n$, explaining the forms shown above. However for these terms care must be taken when any one of $a_1$ through to $a_k$ equals '1'. This is where the map $\theta \otimes \theta$ comes into play. In such cases we symbolically have $0 = a_k - 1 = \delta^{-1}(a_k) \not\in \mathcal{C}$. If $a_1 = 1$ or $a_n = 1$, then the second and penultimate terms in the re-written form above should not be present. If $a_2 = 1$ then the third and fourth terms in the re-written form above should each collapse to $a_1 \otimes a_3 \cdots a_n$. We apply a similar procedure if $a_3 = 1$ or $a_4 = 1$ and so forth up until the case $a_{n-1} = 1$. The action of $\theta \otimes \theta$ precisely enforces these collapses in the special cases mentioned. \hfill \Box

Example 3 Consider the following examples illustrating the co-product $\Delta$ applied to some compositions. Naturally for any integer $n \in \mathcal{C}$ we have $\Delta(n) = \emptyset$, while for other compositions we have, for example, $\Delta(21) = \nu \otimes 21 + 1 \otimes 1 + 21 \otimes \nu$ and also $\Delta(32) = \nu \otimes 32 + 2 \otimes 2 + 3 \otimes 1 + 32 \otimes \nu$, while:

- $\Delta(111) = \nu \otimes 111 + 1 \otimes 1 + 111 \otimes \nu$,
- $\Delta(121) = \nu \otimes 121 + 1 \otimes 11 + 11 \otimes 1 + 121 \otimes \nu$,
- $\Delta(112) = \nu \otimes 112 + 1 \otimes 2 + 11 \otimes 1 + 112 \otimes \nu$,
- $\Delta(241) = \nu \otimes 241 + 1 \otimes 41 + 2 \otimes 31 + 23 \otimes 1 + 241 \otimes \nu$.

The following result establishes that the signature co-algebra is indeed a co-algebra; the proof is given in Appendix $\mathcal{A}$.

Theorem 3 (Co-algebra structure) The signature co-algebra $\mathcal{S}$ is a co-algebra, in particular the co-unit $\varepsilon$ and co-product $\Delta$ satisfy the following defining axioms on $\mathcal{S}$:

(i) $(\text{id} \otimes \varepsilon) \circ \Delta = (\varepsilon \otimes \text{id}) \circ \Delta$ and (ii) $(\text{id} \otimes \Delta) \circ \Delta = (\Delta \otimes \text{id}) \circ \Delta$.

Remark 8 (The signature character map as a homomorphism) With a slight abuse of notation we denote $\chi(\mathcal{S}) \cong \mathbb{R}$ as the space of corresponding signature character values. In this context, we suppose the character map $\chi$ is a homomorphism so that $\chi(u \otimes v) = \chi(u) \chi(v)$ for any pair $u, v \in \mathcal{S}$. 

6 Hierarchy existence and uniqueness

In the P"oppe algebra $\mathbb{R}\langle C \rangle_*$, we introduced in Definition 8 the linear signature expansions $n \in \mathbb{R}\langle C \rangle_*$ given by:

$$ n := \sum_{w \in C(n)} \chi(w) \cdot w. $$

With the formal set up constructed in Sections 4 and 5, we now tackle the following challenge. We have defined the signature expansions $n$ in $\mathbb{R}\langle C \rangle_*$ as linear expansions in monomials of the form $w \in \mathbb{R}\langle C \rangle_*$, The question is, can we find polynomial expansions for the basic monomial single letters $n \in \mathbb{R}\langle C \rangle_*$, in terms of monomials of the form $a_1 \ast a_2 \ast \cdots \ast a_k$? A positive answer guarantees the integrability of the corresponding equation of order $n$ in the non-commutative potential Korteweg-de Vries hierarchy.

Example 4 (Korteweg-de Vries integrability) From their definition, the first three signature expansions are $1 = \chi(1) \cdot 1$, $2 = \chi(2) \cdot 2 + \chi(11) \cdot 11$ and then

$$ 3 = \chi(3) \cdot 3 + \chi(21) \cdot 21 + \chi(12) \cdot 12 + \chi(111) \cdot 111 \iff 3 = 3 \cdot (21 + 12) - 6 \cdot 111. $$

However we also observe that $1 \ast 1 = (\chi(1) \cdot 1) \ast (\chi(1) \cdot 1) = \chi(1) \chi(1) \cdot (1 \ast 1)$ and so

$$ 1 \ast 1 = \chi(1) \chi(1) \cdot (21 + 12 + \chi(q) \cdot (111)) = 21 + 12 + \chi(q) \cdot 111. $$

Substituting the latter result into the former we observe

$$ 3 = 3 - c_{11} \cdot (1 \ast 1), $$

where $c_{11} = 3$. This demonstrates that indeed $3 \in \mathbb{R}\langle C \rangle_*$ has a polynomial expansion in terms of the signature expansions. If we translate this statement back into the Hankel kernel algebra $\mathbb{R}\langle [P_U] \rangle$ it becomes $[U P_3 U] = [U_3 - 3 U_1 U_1]$. This statement establishes that the potential Korteweg-de Vries equations as an integrable Grassmannian flow.

Example 5 (Korteweg-de Vries integrability: fifth order) As in the last example, the question we need to answer is, can the single monomial letter $5 \in \mathbb{R}\langle C \rangle_*$ be expressed as a linear combination of monomials of the form $a_1 \ast a_2 \ast \cdots \ast a_k$? A proposed polynomial of such signature expansions is $\pi_5 = \pi_5 (1, 2, 3, 5)$, where

$$ \pi_5 := c_5 \cdot 5 + c_{11} \cdot 3 \ast 1 + c_{22} \cdot 2 \ast 2 + c_{13} \cdot 1 \ast 3 + c_{111} \cdot 1 \ast 1 \ast 1, $$

where $c_5, c_{11}, c_{22}, c_{13}$ and $c_{111}$ are constants. So the question is, in $\mathbb{R}\langle C \rangle_*$, can we find values for these coefficients $c_5, c_{11}, c_{22}, c_{13}$ and $c_{111}$ such that $5 = \pi_5$? We observe that, given the leading term in $\pi_5$, we should anticipate that $c_5 = 1$. The remaining combination of the quadratic terms $3 \ast 1, 2 \ast 2, 1 \ast 3$ and the cubic term $1 \ast 1 \ast 1$ are chosen because these are the only monomials of this form which when we substitute for the corresponding signature expansions for $1, 2$ and $3$ into them, and compute the P"oppe product of the resulting expansions, we will get compositions in $\mathcal{C}(5)$. And the compositions of $\mathcal{C}(5)$ represent the basis in $\mathbb{R}\langle C \rangle_*$ we should use, guided by the linear signature expansion for $5$ which we can write in the form $5 = 5 - \text{‘lower multi-part compositions’}$. For example, the quadratic term $3 \ast 1$ equals

$$ 3 \ast 1 = (\chi(3) \cdot 3 + \chi(21) \cdot 21 + \chi(12) \cdot 12 + \chi(111) \cdot 111) \ast (\chi(1) \cdot 1) $$

$$ = \chi(3 \circ 1) \cdot 3 \ast 1 + \chi(21 \circ 1) \cdot 21 \ast 1 + \chi(12 \circ 1) \cdot 12 \ast 1 + \chi(111 \circ 1) \cdot 111 \ast 1 $$
\[
\begin{align*}
\chi(3\otimes1) \cdot (41 + 32 + \chi(q) \cdot 311) + \chi(21\otimes1) \cdot (221 + 212 + \chi(q) \cdot 2111) \\
+ \chi(12\otimes1) \cdot (131 + 122 + \chi(q) \cdot 1211) + \chi(111\otimes1) \cdot (1121 + 1112) \\
+ \chi(q) \cdot 11111
\end{align*}
\]

The Pöppe product determines the quadratic and cubic terms in \(\pi_5\). The quadratic terms contain the Pöppe product of the natural numbers shown and all their compositions, to which for the first term in the resulting product a ‘1’ is added to the last letter on the left factor, while for second term a ‘1’ is added to the first letter of the second factor, and for the third term, a letter 1 is squeezed between the two factors. For the cubic term, this process happens twice, for both Pöppe products present, and terms contain the Pöppe product of the natural numbers shown and all their compositions that do not appear on the right in the expression for \(5 \in \mathbb{R}(\mathcal{C})_\pi\) shown at the beginning of this example. In Table 1 we show the coefficients of all the compositions appearing in linear, quadratic and cubic monomials of signature expansions that are in \(\pi_5\). Hence, referring to Table 1 the first column lists all the compositions of 5. The second column shows the signature coefficients of the signature expansion corresponding to 5. The third column shows the signature coefficients of the signature expansion of the polynomial \(3 \ast 1\), once all Pöppe products have been expanded. The fourth and fifth columns show the corresponding signature coefficients in the respective signature expansions of the polynomials \(2 \ast 2\) and \(1 \ast 3\). The sixth column shows the signature coefficients of the signature expansion of \(1 \ast 1 \ast 1\). The seventh (and last) column shows the right-hand side in the equation we are considering, namely \(\pi_5 = 5\). So the question is, can we find the vector of coefficients \(C = (c_5, c_{31}, c_{22}, c_{13}, c_{111})^T\) such that \(\pi_5 = 5\), i.e. such that:

\[
c_5 \cdot 5 + c_{31} \cdot 3 \ast 1 + c_{22} \cdot 2 \ast 2 + c_{13} \cdot 1 \ast 3 + c_{111} \cdot 1 \ast 1 \ast 1 = 5 \quad \iff \quad AC = B,
\]

where \(B\) is the 16 component vector \(B = (\chi(5), 0, \ldots, 0)^T\)? The linear algebraic equation on the right results from equating the coefficients of all 16 compositions appearing in the equation on the left. In the linear algebraic equation, the columns of the \(16 \times 5\) matrix \(A\) are the second through sixth columns of the \(\chi\)-evaluated signatures shown in Table 1. For example the second column of \(A\) is

\[
(\chi(3\otimes1), \chi(3\otimes1), 0, 0, \chi(3\otimes q\otimes1), \chi(21\otimes1), \chi(21\otimes1), \chi(12\otimes1), \ldots)^T.
\]

In the linear equation, columns one through five of \(A\) are naturally associated with the respective coefficients in rows one to five of \(C\).

Let us now outline the strategy we use to solve the linear algebraic equation \(AC = B\). This linear system is overdetermined, there are 16 equations and the components of \(C\) represent the 5 unknowns. Looking at Table 1 it makes sense to swap rows 6 and 7. Let us call the new coefficient matrix incorporating this swap \(A'\). Note the vector \(B\), whose only non-zero component is its first entry, is unaffected by the row swap suggested. For the moment, we ignore the first equation corresponding to the composition \(5 \in \mathbb{R}(\mathcal{C})_\pi\), i.e. we ignore the first row in the augmented matrix \([A' B]\). We focus on the remaining system of 15 homogeneous equations represented by the augmented matrix \([A'(2, \ldots, 16) O]\), where \(A'(2, \ldots, 16)\) represents the submatrix of
Table 1 Non-zero signature coefficients appearing in the expansion of the Pöppe polynomial \( \pi_5 \); see Example 5. The coefficients are the \( \chi \)-images of the signature entries shown. Each column shows the factor contributions to the real coefficients of the compositions of \( C(5) \) shown in the very left column, for each of the monomials in \( \pi_5 \) shown across the top row. The final column represents the right-hand side of the equation \( \pi_5 = 5 \).

| \( C \) | \( 5 \) | \( 3+1 \) | \( 2+2 \) | \( 1+3 \) | \( 1+1+1 \) | \( B \) |
|-------|--------|--------|--------|--------|--------|--------|
| 5     | 5      |        |        |        |        | 5      |
| 41    | 41     | 3\( \otimes 1 \) | 0      |        |        | 0      |
| 32    | 32     | 3\( \otimes 1 \) | 2\( \otimes 2 \) | 0      |        | 0      |
| 23    | 23     | 2\( \otimes 2 \) | 1\( \otimes 3 \) | 0      |        | 0      |
| 14    | 14     |        |        | 1\( \otimes 3 \) | 0      | 0      |
| 311   | 311    | 3\( \otimes q \otimes 1 \) | 2\( \otimes 11 \) | 0      |        | 0      |
| 221   | 221    | 2\( \otimes 1 \) | 2\( \otimes 11 \) | 1\( \otimes 21 \) | 1\( \otimes q \otimes 1 \) | 0      |
| 212   | 212    | 2\( \otimes q \otimes 2 \) | 1\( \otimes 12 \) | 1\( \otimes q \otimes 1 \) | 0      | 0      |
| 131   | 131    | 12\( \otimes 1 \) | 1\( \otimes 12 \) | 1\( \otimes q \otimes 1 \) | 0      | 0      |
| 122   | 122    | 12\( \otimes 1 \) | 1\( \otimes 2 \otimes 2 \) | 1\( \otimes 12 \) | 1\( \otimes q \otimes 1 \) | 0      |
| 113   | 113    |        |        | 1\( \otimes 1 \) | 1\( \otimes 1 \) | 1\( \otimes q \otimes 3 \) | 0      |
| 2111  | 2111   | 2\( \otimes q \otimes 1 \) | 2\( \otimes q \otimes 11 \) | 1\( \otimes 111 \) | 1\( \otimes q \otimes 1 \) | 0      |
| 1211  | 1211   | 12\( \otimes q \otimes 1 \) | 1\( \otimes 111 \) | 1\( \otimes 111 \) | 1\( \otimes q \otimes 1 \) | 0      |
| 1121  | 1121   | 11\( \otimes q \otimes 1 \) | 1\( \otimes 111 \) | 1\( \otimes q \otimes 2 \) | 1\( \otimes q \otimes 1 \) | 0      |
| 1112  | 1112   | 11\( \otimes q \otimes 2 \) | 1\( \otimes q \otimes 2 \) | 1\( \otimes q \otimes 1 \) | 0      | 0      |
| 1111  | 1111   | 11\( \otimes q \otimes 1 \) | 1\( \otimes 111 \) | 1\( \otimes 111 \) | 1\( \otimes q \otimes 1 \) | 0      |

\( A' \) only containing the rows 2 through to 16 and ‘\( O \)’ is a column vector of 15 zeros. For any matrix, its row rank equals its column rank. Further, the rank of \( A'(2, \ldots, 16) \) and \( [A'(2, \ldots, 16) O] \) are the same since the last column of the latter is \( O \). Hence for the linear system of homogeneous equations represented by the augmented matrix \( [A'(2, \ldots, 16) O] \), all but 5 equations are redundant. Indeed let us identify, a-priori, 5 natural equations in \( [A'(2, \ldots, 16) O] \), namely those represented by the first five rows of \( [A'(2, \ldots, 16) O] \); these correspond to the rows in the original matrix \( A \) identified by the compositions 41, 32, 23, 14 and 221. Since the leading diagonal entries in \( A'(2, \ldots, 16) \), i.e. the diagonal entries in \( A'(2, \ldots, 6) \), are all non-zero they can be used as pivots in a Gaussian elimination procedure to render the submatrix \( A'(7, \ldots, 16) \) to the zero submatrix. In other words we can use Gaussian elimination on rows 7 through 16 to ensure \( A'(7, \ldots, 16) = O_{10 \times 10} \). The remaining homogeneous linear system represented by the augmented matrix \( [A'(2, \ldots, 6) O] \) is given by

\[
\begin{pmatrix}
\chi(41) & \chi(3\otimes 1) & 0 & 0 & 0 \\
\chi(32) & \chi(3\otimes 1) & \chi(2\otimes 2) & 0 & 0 \\
\chi(23) & 0 & \chi(2\otimes 2) & \chi(1\otimes 3) & 0 \\
\chi(14) & 0 & 0 & \chi(1\otimes 3) & 0 \\
\chi(221) & \chi(21\otimes 1) & \chi(2\otimes 11) & \chi(1\otimes 21) & \chi(1\otimes 1\otimes 1) \\
\end{pmatrix}
\begin{pmatrix}
c_5 \\
c_{31} \\
c_{22} \\
c_{31} \\
c_{111} \\
\end{pmatrix} = \begin{pmatrix} 0 \\ 0 \\ 0 \\ 0 \\ 0 \end{pmatrix}.
\]

Using Definition 5 and that \( \chi \) is a homomorphic map we have \( \chi(41) = \chi(14) = 5 \), \( \chi(32) = \chi(23) = 10 \), \( \chi(221) = 30 \), \( \chi(3\otimes 1) = \chi(3)\chi(1) = 1 \), \( \chi(2\otimes 1) = \chi(2)\chi(1) = 3 \), and so forth. Such a homogeneous linear system has a solution. The question is whether the solution is the unique trivial solution \( C = O \), or there is a general solution with one or more free variables. See Meyer [51, p. 61] for more details. We observe the final column in \( A'(2, \ldots, 6) \) just above is a basic column that allows us to determine the final variable \( c_{111} \) in terms of the other unknowns \( c_5, c_{31}, c_{22} \) and \( c_{13} \). Hence we focus
on the $4 \times 4$ subsystem for the latter four variables. By Gaussian elimination the $4 \times 4$ homogeneous subsystem reduces to $[A^\prime O]$ where $A^\prime$ is given by

$$
\begin{pmatrix}
5 & 1 & 0 & 0 \\
0 & 1 & -1 & 0 \\
0 & 0 & 1 & -1 \\
0 & 0 & 0 & 0
\end{pmatrix},
$$

and there is thus one free variable. The other three variables are given explicitly, linearly, in terms of that free variable. For example we can express the solution as $c_{31} = -5c_5$, $c_{22} = -5c_4$ and $c_{13} = -5c_3$. The free variable is fixed by the very first equation in the original linear system $AC = B$ which is $c_5 = 1$. Since the final equation in the linear system $[A^\prime(2, \ldots , 6)O]$ giving $c_{111}$ is $c_{111} = -30c_5 - 3c_{11} - 2c_{22} - 3c_{13}$, we find $c_{111} = 10$. Hence we have solved the linear system of equations $A^\prime C = B$. The unique solution is $C = (1, -5, -5, -5, 10)^T$. Hence we have shown that,

$$\pi_5 := 5 - 5 \cdot (3 \ast 1 + 2 \ast 2 + 1 \ast 3) + 10 \cdot (1 \ast 1 \ast 1).$$

We now turn to our main result. Example 5 illustrated the strategy we employ for the general order case.

**Definition 12 (Pöppe polynomials)** For $n \in \mathbb{N}$, let $\pi_n = \pi_n(1, 2, \ldots, (n - 2), n)$ denote a general polynomial consisting of a linear combination of monomials of signature expansions of the following form:

$$\pi_n := \frac{1}{2}(n+1) \sum_{a_1a_2 \cdots a_k \in \mathcal{C}^\prime(n)} c_{a_1a_2 \cdots a_k} \ast a_1 \ast a_2 \ast \cdots \ast a_k,$$

where $\mathcal{C}^\prime(n) \subset \mathcal{C}(n)$ represents the subset of compositions $w = a_1a_2 \cdots a_k$ of $n$ such that $a_1 + a_2 + \cdots + a_k = n - k + 1$. The coefficients $c_{a_1a_2 \cdots a_k}$ are real constants. We call polynomials of this form in $\mathbb{R}(\mathcal{C})_+$, Pöppe polynomials.

Note, in the definition of the Pöppe polynomials, $\mathcal{C}^\prime(n)$ is the correct subset of $\mathcal{C}(n)$ for the polynomial expansion $\pi_n$, as we will equate $\pi_n$ to compositions of $n$ and each Pöppe product $\ast$ adds 1 to the eventual compositions of $n$ produced in the products $a_1 \ast a_2 \ast \cdots \ast a_k$. The product $a_1 \ast a_2 \ast \cdots \ast a_k$ contains $k - 1$ such Pöppe products. After expanding each of the Pöppe polynomial factors and expanding all the $k - 1$ Pöppe products between the resulting terms, to guarantee all the compositions in the final sum are indeed compositions of $n$, we must restrict the sum of the composition digits $a_1 + a_2 + \cdots + a_k$ to be $n - (k - 1)$. Note the upper bound for $k$ is $\frac{1}{2}(n+1)$ as this corresponds to the $k$-length monomial $1 \ast 1 \ast \cdots \ast 1$ generating compositions of $n$.

**Theorem 4 (Main result: Integrability)** For every odd natural number $n$, there exists a unique set of real coefficients $\{c_w : w \in \mathcal{C}^\prime(n)\}$ such that in $\mathbb{R}(\mathcal{C})_+$, we have

$$n = \pi_n(1, 2, \ldots, (n - 2), n).$$

Before proving this theorem, we introduce a 'descent ordering' on compositions as well as some useful notation.
We also use the notation $C_n^2$ which start with the digit ‘(n, k)’ to denote the set of $k$-part compositions of $n$, i.e. the subset of compositions $C(n)$ whose length equals $k$: $w \in C(n, k)$ if $w \in C(n)$ and $w = a_1 a_2 \cdots a_k$. We note, naturally we have $C(n) = \cup_{k=1}^{n} C(n, k)$ and so forth. Here the notation $(n, k)$, say $k \in \mathbb{N}$, then $u$ precedes $v$ if for some $\ell \in \{1, 2, \ldots, k\}$ we have $u_1 = v_1, u_2 = v_2, \ldots, u_{\ell-1} = v_{\ell-1}$ and $u_\ell < v_\ell$. Otherwise we have $v \prec u$. The resulting ordering induced on $C$, is the descent ordering.

The descent ordering on $C(n)$ is naturally transferred to subsets of $C(n)$ such as $C^\ast(n)$. We also use the notation $C(n, k)$ to denote the set of $k$-part compositions of $n$, i.e. the subset of compositions $C(n)$ whose length equals $k$: $w \in C(n, k)$ if $w \in C(n)$ and $w = a_1 a_2 \cdots a_k$. We note, naturally we have $C(n) = \cup_{k=1}^{n} C(n, k)$ and so forth. Here the notation $(n, k)$, say $k \in \mathbb{N}$, then $u$ precedes $v$ if for some $\ell \in \{1, 2, \ldots, k\}$ we have $u_1 = v_1, u_2 = v_2, \ldots, u_{\ell-1} = v_{\ell-1}$ and $u_\ell < v_\ell$. Otherwise we have $v \prec u$. The resulting ordering induced on $C$, is the descent ordering.

**Definition 13 (Descent ordering of compositions)** A composition $u \in C$ precedes another composition $v \in C$ and we write ‘$u \prec v$’, if the length of the composition $u$, i.e. the number of digits it contains, is strictly less than the length of $v$. If $u$ and $v$ have the same length, say $k$, so $u = u_1 u_2 \cdots u_k$ and $v = v_1 v_2 \cdots v_k$, then $u$ precedes $v$ if for some $\ell \in \{1, 2, \ldots, k\}$ we have $u_1 = v_1, u_2 = v_2, \ldots, u_{\ell-1} = v_{\ell-1}$ and $u_\ell < v_\ell$. Otherwise we have $v \prec u$. The resulting ordering induced on $C$, is the descent ordering.

The descent ordering on $C(n)$ is naturally transferred to subsets of $C(n)$ such as $C^\ast(n)$. We also use the notation $C(n, k)$ to denote the set of $k$-part compositions of $n$, i.e. the subset of compositions $C(n)$ whose length equals $k$: $w \in C(n, k)$ if $w \in C(n)$ and $w = a_1 a_2 \cdots a_k$. We note, naturally we have $C(n) = \cup_{k=1}^{n} C(n, k)$ and so forth. Here the notation $(n, k)$, say $k \in \mathbb{N}$, then $u$ precedes $v$ if for some $\ell \in \{1, 2, \ldots, k\}$ we have $u_1 = v_1, u_2 = v_2, \ldots, u_{\ell-1} = v_{\ell-1}$ and $u_\ell < v_\ell$. Otherwise we have $v \prec u$. The resulting ordering induced on $C$, is the descent ordering.

We are now in a position to prove Theorem 4.

**Proof** We use Example 5 and Table 1 as an example template for our strategy. Imagine we construct the array in Table 1 for the general integer case $n \in \mathbb{N}$ rather than the $n = 5$ case shown in the table. We parametrise the rows of the table using the set of compositions $C(n)$, ordering the compositions according to the
descent ordering given in Definition 13. This naturally splits the rows into blocks of compositions of the same length (i.e. number of parts) according to the decomposition,

\[ C(n) = C(n, 1) \cup C(n, 2) \cup \cdots \cup C(n, n). \]

This is the highest level coarse-grained block decomposition of \( C(n) \). We parametrise the columns of the table using the set of decompositions \( C^*(n) \). We again use descent ordering to parametrise the columns, and this naturally splits the columns into blocks of compositions of \( C^*(n) \) with the same number of parts according to the decomposition,

\[ C^*(n) = C^*(n, 1) \cup C^*(n - 1, 2) \cup C^*(n - 2, 3) \cup \cdots \cup C^*\left(\frac{n}{2}(n + 1), \frac{1}{2}(n + 1)\right). \]

Recall our goal is to prove there exists a unique set of real coefficients \( \{c_w : w \in C^*(n)\} \) such that in \( R\langle C^* \rangle \), we have

\[ \sum_{k=1}^{\frac{n}{2}(n+1)} \sum_{a_1, a_2, \ldots, a_k \in C^*(n, k)} c_{a_1a_2\ldots a_k} \cdot a_1 \ast a_2 \ast \cdots \ast a_k = n. \]

Each column of the overall table corresponds to a term \( a_1 \ast a_2 \ast \cdots \ast a_k \) and its associated coefficient \( c_{a_1a_2\ldots a_k} \) from the sum on the left. The final column in the overall table represents the right-hand side in the equation above which is merely the single letter composition \( n \in R\langle C^* \rangle \). Each term \( a_1 \ast a_2 \ast \cdots \ast a_k \), once the Pöppe products have been evaluated, generates a linear combination of compositions from \( C(n) \). Within each column corresponding to \( a_1a_2\cdots a_k \in C^*(n) \), each row contains the coefficient of the corresponding composition generated in the evaluation of the product \( a_1 \ast a_2 \ast \cdots \ast a_k \).

In other words, given a composition representing a given row in the table, the element of that row in the column associated with a composition \( a_1a_2\cdots a_k \in C^*(n) \), contains the \( \chi \)-image of the corresponding signature of the composition corresponding to that row. Hence the resulting overall table, just like Table 11 represents a linear equation of the form

\[ AC = B, \]

where \( C \) is the vector of unknown coefficients \( \{c_w : w \in C^*(n)\} \), in descent order so \( C = (c_n, c_{(n-1)}, c_{(n-2)}, \ldots, c_{1}, c_{(n-2)}1, \ldots) \). Since the number of \( k \)-part compositions of \( n \) is \( (n-1) \) choose \((k-1)\), the length of the vector \( C \) is

\[ |C| = \sum_{k=1}^{\frac{n}{2}(n+1)} \binom{n - (k - 1)}{k}. \]

The vector \( B \) contains all zeros apart from the very first element which is the \( \chi \)-image of the composition element \( n \in C(n) \), i.e. \( B = (\chi(n), 0, 0, \ldots)^T \). Naturally we have \( \chi(n) = 1 \). The number of rows in \( A \) and the length of \( B \) is \( 2^{n-1} \), the total number of compositions of \( n \). Hence \( A \) is a matrix of size \( 2^{n-1} \times |C| \), with \( |C| < 2^{n-1} \) for \( n \geq 2 \).

As already indicated, we can decompose rows of \( A \) into blocks of compositions in \( C(n,k') \) for \( k' = 1, \ldots, n \), and we can decompose columns of \( A \) into blocks of compositions in \( C^*(n-k+1,k) \) for \( k = 1, \ldots, \frac{n}{2}+1 \). We note the term \( a_1 \ast a_2 \ast \cdots \ast a_k \) with \( a_1a_2\cdots a_k \in C^*(n-k+1,k) \), once the Pöppe products have been evaluated, only generates compositions in \( C(n) \) with \( k \) or more parts. Thus in terms of these blocks, parametrised by \((k',k)\), the matrix \( A \) is lower block triangular. Let us now
where we show the non-zero elements in this block. Note for the rows parametrised by the columns parametrised by the rows in the augmented system \( A B \), we systematically determine the coefficients \( c_{a_1 a_2 \cdots a_k} \) corresponding to each block.

First consider the case \( k = 1 \), corresponding to the column block \( C^\ast(n, 1) \) and thus the element \( n \) and coefficient \( c_n \). Since we know \( n = \sum_{w \in C(n)} \chi(w) \cdot w \), this means the first row of this single column block contains \( \chi(n) = 1 \), while the subsequent rows corresponding to \( w \in C(n) \) contain the coefficient \( \chi(w) \), in descent order. The first row of the linear system represented by the augmented matrix \( [A B] \) represents the equation \( c_n = 1 \). However as we did in Example \ref{example} we temporarily ignore this knowledge.

Second we focus on the 2-part composition sub-block of \( A \) parametrised by the columns \( C^\ast(n, 1) \cup C^\ast(n-1, 2) \) and by the rows \( C(n, 2) \), both in descent order. See Table \ref{table} where we show the non-zero elements in this block. Note for the rows parametrised by \( C(n, 2) \) all the columns in \( A \) to the right of the block contain zero entries and the corresponding entries in \( B \) are all zero. This block has size \((n-1) \times (n-1)\). If we focus on these rows in the augmented system \([A B]\) then we have \((n-1)\) equations in the \((n-1)\) unknowns \( c_n, c_{(n-2)1}, c_{(n-3)2}, \ldots, c_{1(n-2)} \). We also observe that all the leading diagonal entries in the block shown in Table \ref{table} from the top left \((1, 1)\) entry to the lower right \((n-1, n-1)\) entry, are all non-zero. Thus all these leading diagonal entries can be used as pivots in a Gaussian elimination process to render all the entries in the corresponding columns in the rows below this block to be zero. We can also solve the homogeneous linear system above, represented by this block, to find the coefficients.

For \( \ell = 2, \ldots, (n-2) \) the homogeneous linear system of equations takes the form,

\[
\begin{align*}
\chi((n-1)1) c_n + \chi((n-2)\otimes 1) c_{(n-2)1} &= 0, \\
\chi((n-\ell)\ell) c_n + \chi((n-\ell)\otimes (\ell-1)) c_{(n-\ell)(\ell-1)} + \chi((n-\ell-1)\otimes \ell) c_{(n-\ell-1)\ell} &= 0, \\
\chi(1(n-1)) c_n + \chi(1\otimes (n-2)) c_{1(n-2)} &= 0.
\end{align*}
\]

Note, apart from the coefficients of \( c_n \) shown, all the other coefficients of \( c_{(n-2)1}, c_{(n-3)2}, \ldots, c_{1(n-2)} \) are equal to one. Starting with the first equation above, we can iteratively solve for all the coefficients in terms of \( c_n \) giving, for \( \ell = 1, \ldots, (n-2) \), the relations \( \hat{c}_\ell c_n + c_{(n-\ell-1)\ell} = 0 \), where

\[
\hat{c}_\ell := \chi((n-\ell)\ell) - \chi((n-\ell+1)(\ell-1)) + \chi((n-\ell+2)(\ell-2)) - \cdots + (-1)^{\ell+1} \chi((n-1)1),
\]

along with the final equation which is \( \chi(1(n-1)) c_n + c_{1(n-2)} = 0 \). We now have two linear homogeneous equations relating the unknowns \( c_{1(n-2)} \) and \( c_n \), the equation for which \( \ell = n-2 \) and the final equation just mentioned. These two equations are

| \( n \) | \( (n-2) \times 1 \) | \( (n-3) \times 2 \) | \( (n-4) \times 3 \) | \( \cdots \) | \( 2 \times (n-3) \) | \( 1 \times (n-2) \) |
|------|----------|----------|----------|----------|----------|----------|
| \( (n-1)1 \) | \( (n-2)\otimes 1 \) | \( \cdots \) | \( \cdots \) | \( \cdots \) | \( 2\otimes(n-3) \) | \( 1\otimes(n-2) \) |
| \( (n-2)2 \) | \( (n-2)\otimes 1 \) | \( (n-3)\otimes 2 \) | \( \cdots \) | \( \cdots \) | \( \cdots \) | \( \cdots \) |
| \( (n-3)3 \) | \( (n-3)\otimes 2 \) | \( (n-4)\otimes 3 \) | \( \cdots \) | \( \cdots \) | \( \cdots \) | \( \cdots \) |
| \( \vdots \) | \( \vdots \) | \( \vdots \) | \( \vdots \) | \( \vdots \) | \( \vdots \) | \( \vdots \) |
| \( 2(n-2) \) | \( \cdots \) | \( 2\otimes(n-3) \) | \( 1\otimes(n-2) \) | \( \cdots \) | \( \cdots \) | \( \cdots \) |
| \( 1(n-1) \) | \( \cdots \) | \( 1\otimes(n-2) \) | \( \cdots \) | \( \cdots \) | \( \cdots \) | \( \cdots \) |

Table 2 Non-zero signature coefficients appearing in the 1-part and 2-part composition blocks of \( A \), ignoring the first row. The coefficients are the \( \chi \)-images of the signature entries shown.
consistent if $\chi_{n-2} = \chi(1(n-1))$. Note from its definition, $\chi(nk)$ is the binomial coefficient $n+k$ choose $n$, and thus equal to $\chi(kn)$. Depending on whether $\ell$ is odd or even we have,

$$\hat{X}_\ell = \begin{cases} 2m + 1, & \text{if } \ell = 2m + 1, \\ 2m - 2, & \text{if } \ell = 2m. \end{cases}$$

To see this, we observe that in the odd case we can utilise the symmetry $\chi(nk) = \chi(kn)$ and the differing signs in front of each as they appear in $\hat{X}_\ell$, to cancel all the terms except the final one given by $(-1)^{2m}\chi((2m)1) = 2m + 1$. In the even case we observe

$$\hat{X}_{2m} = \left( \begin{array}{c} 2m \\ \frac{2m}{3} \\ \frac{2m}{4} \\ \vdots \\ \frac{2m}{2m-1} \end{array} \right)$$

$$= \left( \frac{1 - x}{x} \right)^{2m - 1} + (2m)x - x^{2m} \bigg|_{x=1}$$

$$= 2m - 2.$$

Hence we observe when $n-2 = 2m+1$ we have $\hat{X}_{n-2} = \hat{X}_{2m+1} = 2m+1 = \chi(1(n-1))$ and the homogeneous linear system of equations above is consistent with a single free variable, namely $c_{n}$. On the other hand when $n-2 = 2m$ then $\hat{X}_{2m} = 2m - 2$ and this does not equal $\chi(1(n-1)) = 2m$. Thus in the latter case there is no solution to the linear homogeneous system $[AB]$ corresponding to the 2-part composition rows. Given the consistency in the odd order case, we now recall the equation from the very first row of the linear system $[AB]$ which states $c_{n} = 1$. When we include this information we observe we can solve the top part of the linear system of equations $[AB]$, corresponding to the single and then 2-part composition rows, uniquely for the coefficients $c_{0}, c_{(n-2)1}, c_{(n-3)2}, \ldots, c_{1(n-2)}$.

Third we now focus on solving the linear system $[AB]$ for all the coefficients $c_{w}$ with $w \in C^*(n)$ systematically. We assume $n$ is odd, and in some sense we start again from the beginning and ignore the information given in the equation represented by the very first row, i.e. that $c_{n} = 1$. We focus on the homogeneous linear system given by the remaining rows and represented by $[A(2, \ldots, 2^{n-1}) O]$, where as before $A(2, \ldots, 2^{n-1})$ represents the submatrix of $A$ consisting of rows 2 through $2^{n-1}$ As is evident from the previous paragraph, the homogeneous linear system of $(n-1)$ equations represented by $[A(2, \ldots, n) O]$ in the $(n-1)$ unknowns $c_{n}, c_{(n-2)1}, c_{(n-3)2}, \ldots, c_{1(n-2)}$ can be rendered into an upper triangular form with non-zero entries on the leading diagonal. Further, in the rows 2 through $n$ considered, the entries in the columns parametrised by $C^*(2-3) \cup C^*(3-4) \cup \cdots \cup C^*(\frac{1}{2}(n+1), \frac{1}{2}(n+1))$ are all zero.

Continuing, we now focus on the 3-part composition sub-block of $A$ parametrised by the columns $C^*(n-2, 3)$ and by the rows $C(n, 3)$, both in descent order. The reader may find Table 3 in Appendix B helpful in visualising the block structure we now discuss. Note all the entries in $A$ above and to the right of this sub-block are zero. Our goal here is to show, the 3-part composition rectangular sub-block of $A$ is itself lower triangular, in the sense that by row swaps we rearrange the sub-block so that each column has a non-zero pivot entry on the leading diagonal with all the entries in the rows above the pivot equal to zero. If this is possible, then we can use elementary row operations to render the entries in all the rows below leading diagonal rows, corresponding to any compositions with 3- or more parts, to be zero. Note we can use elementary row operations to render any such redundant rows below the 2-part composition rows in the columns parametrised by the 1- and 2-part composition columns,
i.e. by $C^*(n, 1) \cup C^*(n - 1, 2)$, to have zero entries. This is because the leading diagonal in the homogeneous linear system corresponding to these columns has non-zero entries. Thus using the leading diagonal rows in the 3-part composition sub-block of $A$, we can uniquely determine the unknowns $c_w$ with $w \in C^*(n - 2, 3)$ in terms of $c_{n}$, just as we achieved (see above) for the unknowns $c_w$ with $w \in C^*(n - 1, 2)$. Recall, retaining descent order, we express

$$C(n, 3) = \bigcup_{k' = 1}^{n-2} (n - k' - 1)C(k' + 1, 2) \quad \text{and} \quad C^*(n - 2, 3) = \bigcup_{k = 1}^{n-4} (n - k - 3) \ast C^*(k + 1, 2).$$

These two decompositions of $C(n, 3)$ and $C^*(n - 2, 3)$ show we can decompose the corresponding sub-block of $A$ into further, level 2 sub-blocks. The level 2 sub-blocks are parametrised by $k'$ and $k$ with the ranges indicated above. Further, each level 2 sub-block, for each $(k', k)$ pair, is, within themselves parametrised by the 2-part composition rows and columns $C(k' + 1, 2)$ and $C^*(k + 1, 2)$, respectively. We proceed systematically, pairing up level 2 sub-blocks. To begin, we observe the top left level 2 sub-block which corresponds to the values $(k', k) = (1, 1)$ and is a $1 \times 1$ level 2 sub-block corresponding to the row $(n - 2)11$ and column $(n - 4)11$ must have a zero entry as $(n - 4)111$ cannot generate the composition $(n - 2)11$. Indeed, by the same argument, all the entries in the remaining columns in this row, i.e. to the right in terms of descent order, must also be zero. Hence the homogeneous linear equation represented by this row only contains the unknowns $c_w$ with $w \in C^*(n, 1) \cup C^*(n - 1, 2)$. Since we already have a set of $(n - 1)$ homogeneous linear equations in these $(n - 1)$ unknowns with non-zero diagonal entries, the new equation must be a linear combination of these equations, and thus contains no new information. Alternatively, as discussed above, all the entries in this row in preceding columns corresponding to the 1- and 2-part compositions $C^*(n, 1) \cup C^*(n - 1, 2)$ can be rendered zero by elementary row operations. Thus again the whole row contains no new information and we can ignore it, which we do henceforth. Further since for $k = 1, \ldots, n - 4$ the elements in $(n - k - 3) \ast C^*(k + 1, 2)$ can only generate compositions $(n - k' - 1)C(k' + 1, 2)$ when $(n - k' - 1) - (n - k - 3) \leq 1$, which is equivalent to $k \leq k' - 1$, we deduce that for $k \geq k'$, all the level 2 sub-blocks are zero. Hence in terms of level 2 sub-blocks the 3-part composition rectangular sub-block of $A$ is lower triangular. Now we just need to demonstrate each level 2 sub-block parametrised by $(k', k) = (k + 1, k)$ itself is lower triangular. Each such level 2 sub-block consists of the rows $(n - k - 2)C(k + 2, 2)$ and columns $(n - k - 3) \ast C^*(k + 1, 2)$. The first letters in these compositions are fixed and the rows and columns within the level 2 sub-block are parametrised by $C(k + 2, 2)$ and $C^*(k + 1, 2)$. It is straightforward to check any such level 2 rectangular sub-blocks are indeed lower triangular with non-zero entries on the diagonal as follows. Recalling the structure of the level 1 sub-block of 2-part composition rows and columns from above, we know that each such level 2 rectangular sub-block has size $(k + 1) \times k$. The final row has a non-zero entry in the final column of the sub-block. However as we proceed through $k = 1, 2, 3, \ldots$ we observe that each homogeneous linear equation represented by those final rows must be a linear combination of the homogeneous linear equations that precede it, as those preceding equations are a square system of linear combinations of the same set of unknowns. Hence each such final row is redundant. Hence the 3-part composition rectangular sub-block of $A$ is lower triangular and we can solve for the unknowns $c_w$ with $w \in C^*(n - 2, 3)$ linearly in terms of $c_{n}$. Note that all entries in the columns $C^*(n - 2, 3)$ below the rows $C(n, 3)$ can be rendered zero by elementary row operations.
We now briefly outline the procedure for determining the unknowns $c_w$ with $w \in C^*(n-3,4)$ linearly in terms of $c_n$. As we do so, the nature of the procedure for determining the remaining unknowns becomes apparent, and reveals itself to be recursive and straightforward. We focus on the 4-part composition sub-block of $A$ parametrised by the columns $C^*(n-3,4)$ and by the rows $C(n,4)$, both in descent order. The reader may find Table 4 in Appendix B helpful in visualising the block structure we discuss herein.

All the entries in $A$ above and to the right of this sub-block are zero, since elements in $C^*(n-3,4)$ can only generate compositions in $C(n)$ with four or more parts. Our goal is to show, 4-part composition rectangular sub-block of $A$ is itself lower triangular, as we did for the 3-part composition rectangular sub-block of $A$. As above, if this is possible, then we can use elementary row operations to render the entries in all the rows below the leading diagonal rows, corresponding to any compositions with 4- or more parts, to be zero. The leading diagonal rows in the 4-part composition sub-block of $A$, uniquely determine the unknowns $c_w$ with $w \in C^*(n-3,4)$ in terms of $c_n$. Recall, as always retaining descent order, we have

$$C(n,4) = \bigcup_{k'=1}^{n-3} (n-k'-2)C(k'+2,3) \quad \text{and} \quad C^*(n-3,4) = \bigcup_{k=1}^{n-6} (n-k-5) \ast C^*(k+2,3).$$

Thus we observe we decompose the sub-block of $A$ parametrised by the rows $C(n,4)$ and columns $C^*(n-3,4)$ into level 2 sub-blocks. The level 2 sub-blocks are parametrised by $k'$ and $k$ as indicated above. Each level 2 sub-block, for each $(k',k)$ pair, is, within themselves parametrised by the 3-part composition rows and columns $C(k'+2,3)$ and $C^*(k+2,3)$, respectively. And each of these level 2 sub-blocks can be decomposed into further level 3 sub-blocks of 2-part compositions, as we demonstrated in the previous paragraph. We start by showing that in terms of level 2 sub-blocks, the 4-part composition sub-block of $A$ is lower triangular. Observe the top left level 2 sub-block which corresponds to the values $(k',k) = (1,1)$ is a $1 \times 1$ level 2 sub-block corresponding to the row $(n-3)111$ and column $(n-6)111$ must have a zero entry as $(n-6) \ast 1 \ast 1 \ast 1$ cannot generate the composition $(n-3)111$. Indeed, by the same argument, all the entries in the remaining columns in this row, i.e. to the right in terms of descent order, must also be zero. Hence the homogeneous linear equation represented by this row does not generate any new information. Next, the sub-block which corresponds to the values $(k',k) = (2,1)$ is a $3 \times 1$ level 2 sub-block corresponding to the rows $(n-4)211$, $(n-4)121$ and $(n-4)112$ and column $(n-6)111$ must have zero entries as $(n-6) \ast 1 \ast 1 \ast 1$ cannot generate any of the compositions $(n-4)211$, $(n-4)121$ and $(n-4)112$. By similar arguments these rows do not generate any new information either. Indeed, we observe that for $k = 1, \ldots, n-6$ the elements in $(n-k-5) \ast C^*(k+2,3)$ can only generate compositions $(n-k'-2)C(k'+2,3)$ when $(n-k'-2) - (n-k-5) \leq 1$, which is equivalent to $k \leq k'-2$, we deduce that for $k \geq k'-1$, all the level 2 sub-blocks are zero. Thus, in terms of level 2 sub-blocks, the 4-part composition rectangular sub-block of $A$ is lower triangular—either by deleting the rows corresponding to the first two row sub-blocks or by row-swapping them to rows sufficiently far down the matrix. Now we demonstrate that each level 2 sub-block with $k' = k+2$ is itself lower triangular. Indeed we have already demonstrated this in the last paragraph, which recursively used that each of these level 2, 3-part composition sub-blocks can be further decomposed into level 3, 2-part composition sub-blocks, which are lower triangular. Consequently we can solve for the unknowns $c_w$ with $w \in C^*(n-3,4)$ linearly in terms of $c_n$. 

Finally we discuss the general case, for any \( 5 \leq \ell \leq \frac{1}{2}(n+1) \), for the \( \ell \)-part composition sub-block of \( A \) parametrised by the columns \( C^*(n-\ell+1, \ell) \) and by the rows \( C(n, \ell) \), both in descent order. All the entries in \( A \) above and to the right of this sub-block are zero, since elements in \( C^*(n-\ell, \ell) \) can only generate compositions in \( C(n) \) with \( \ell \)-parts or more. Assume that sub-blocks parametrised by rows \( C(n, \ell-1) \) and columns \( C^*(n-\ell+2, \ell-1) \), i.e. at the previous stage ‘\( \ell-1 \)’, are lower triangular. Now, again retaining descent order, we have

\[
C(n, \ell) = \bigcup_{k'=1}^{n-\ell+1} (n-k'+\ell+2)C(k'+\ell-2, \ell-1)
\]

and

\[
C^*(n-\ell+1, \ell) = \bigcup_{k=1}^{n-2\ell+2} (n-2\ell+3-k)*C^*(k+\ell-2, \ell-1).
\]

Thus we decompose the sub-block of \( A \) parametrised by the rows \( C(n, \ell) \) and columns \( C^*(n-\ell+1, \ell) \) into level 2 sub-blocks. The level 2 sub-blocks are parametrised by \( k' \) and \( k \) as indicated above. For \( k = 1, \ldots, n-2\ell+2 \) the elements in \( (n-2\ell+3-k)*C^*(k+\ell-2, \ell-1) \) can only generate compositions of \( (n-k'-\ell+2)C(k'+\ell-2, 3) \) when \( (n-k'-\ell+2)-(n-2\ell+3-k) \leq 1 \), which is equivalent to \( k \leq k' - \ell + 2 \). Hence for \( k \geq k' - \ell + 3 \), all the level 2 sub-blocks are zero. Thus, in terms of level 2 sub-blocks, the \( \ell \)-part composition rectangular sub-block of \( A \) is lower triangular—or can be transformed as such by elementary row swaps. Now using that we assumed the sub-blocks parametrised by rows \( C(n, \ell-1) \) and columns \( C^*(n-\ell+2, \ell-1) \) are lower triangular, the sub-block of \( A \) parametrised by the columns \( C^*(n-\ell+1, \ell) \) and the rows \( C(n, \ell) \) must be lower triangular. Hence we can solve for the unknowns \( c_w \) with \( w \in C^*(n-\ell+1, \ell) \) linearly in terms of \( c_n \). By induction, this applies for all values of \( \ell \) up to and including \( \ell = \frac{1}{2}(n+1) \). Finally now utilising that \( c_n = 1 \) generates a unique solution set for \( c_w \) for all \( w \in C^*(n) \).

\[\square\]

**Remark 10 (Even order)** For every even natural number \( n \), no such set of real coefficients exist. This is demonstrated in the proof above in the component addressing the 2-part composition sub-blocks of \( A \).

**Remark 11 (Linear combinations of fields)** Since we can expand each of the single letter compositions in \( \mu_3 \cdot 3 + \mu_5 \cdot 5 + \mu_7 \cdot 7 + \cdots \), in terms of \( \mu_3 \cdot \pi_3 + \mu_5 \cdot \pi_5 + \mu_7 \cdot \pi_7 + \cdots \), we deduce that we can solve/linearise, in the manner we have outlined, any nonlinear evolutionary partial differential equation generated by a linear combination of fields from the non-commutative primitive Korteweg-de Vries hierarchy.

### 7 Lax Hierarchy

Our goal in this section is to show how the non-commutative Lax hierarchy is easily and naturally deduced using the signature expansions in the Pöppe algebra \( \mathbb{R}(\mathcal{C}) \) from Sections 4 and 6. We derive the well-known non-commutative Lax hierarchy, see for example Carillo and Schoenlieb [16], and show that at each odd order, the underlying iteration generates nonlinear terms that can be expressed as a linear combination of monomials of the form \( a_1 \cdot a_2 \cdot \cdots \cdot a_k \) with \( a_1 a_2 \cdots a_k \in C^*(n) \). By our main result
Theorem 4 from the last section on the existence and uniqueness of the Korteweg–de Vries hierarchy, we deduce the non-commutative hierarchy we developed therein, and the non-commutative Lax hierarchy, are one and the same.

We proceed by deriving the first part of the result on the non-commutative Lax hierarchy, before stating the main result of this section, and then subsequently proving the second part. For a given odd natural number $n \in \mathbb{N}$, we assume there exists a polynomial $\hat{\pi}_n = \hat{\pi}_n(1, 2, \ldots, (n - 2))$ such that the composition element $n \in \mathbb{R}\langle C \rangle_*$ can be expressed in the form

$$n = n + \hat{\pi}_n(1, 2, \ldots, (n - 2)).$$

For $n = 3$, such a polynomial $\hat{\pi}_3 = \hat{\pi}_3(1)$ exists as we know $3 = 3 - 3 \cdot 1 \cdot 1$, representing the potential Korteweg–de Vries equation, i.e. $\hat{\pi}_3 = -3 \cdot 1 \cdot 1$. We now apply the derivation operator $\partial$ to the expression for the composition element $n \in \mathbb{R}\langle C \rangle_*$ just above, twice. Recall the action of $\partial$ on monomials in $\mathbb{R}\langle C \rangle_*$ from Definition 5. Starting with $\partial(n) = 1n + (n + 1) + n1$ and then $\partial(1n) = 2 \cdot (11n) + 2n + 1(n + 1) + 1n1$ and $\partial(n1) = 1n1 + (n + 1)1 + 2 \cdot (n11) + n2$, the successive application of $\partial$ generates, first that $(n + 1) + (1n + n) = (n + 1) + \partial \hat{\pi}_n$, and then:

$$(n + 2) + 2 \cdot (11n + 1(n + 1) + 1n1 + (n + 1)1 + n11) + (2n + n2) = (n + 2) + \partial^2 \hat{\pi}_n.$$

Next, in the last expression, looking at the terms on the left, not including the first term, and recalling the properties of the Pöppe product, we observe

$$1 \ast n + n \ast 1 = 2n + 1(n + 1) + 2 \cdot (11n) + (n + 1)1 + n2 + 2 \cdot (n11).$$

Substituting this result into the previous expression and rearranging, we find,

$$(n + 2) = (n + 2) + \partial^2 \hat{\pi}_n - (1 \ast n + n \ast 1) - (1(n + 1) + (n + 1)1 + 2 \cdot (n11)).$$

To find a closed expression for the terms '$1(n + 1) + (n + 1)1 + 2 \cdot (n11)' in terms of a Pöppe polynomial, we use the following result relating the Pöppe product and the derivation operator $\partial$.

**Lemma 6** Given two arbitrary compositions $ua, bv \in \mathbb{R}\langle C \rangle_*$ where we distinguish the final and beginning single letters $a$ and $b$ as shown, we have:

$$ua \ast bv = \partial(ubav) - (\partial u)abv - uab(\partial v) + uabv.$$

**Proof** From the definition of $\partial$ we have $\partial(uabv) = (\partial u)abv + u(\partial(ab))v + uab(\partial v)$. Thus from the definition of the Pöppe product, the result follows as we observe $ua \ast bv$ equals $u(a + 1)bv + ua(b + 1)v + 2 \cdot (ua1bv)$, which in turn can be expressed as $u(\partial(ab))v + uabv$, which itself in turn can be expressed as $\partial(ubav) - (\partial u)abv - uab(\partial v) + uabv$. $\square$

Now we observe that since $\partial(n) = 1n + (n + 1) + n1$ then utilising Lemma 6 with $ua = 1$ and then $bv$ successively equal to $1n$, $(n + 1)$ and $n1$, and vice-versa, we have:

$$1 \ast \partial(n) + \partial(n) \ast 1 = \partial(1(n + 1) + (n + 1)1 + 2 \cdot (1n1) + 1n + n11) - \left(11d(n) + \partial(n)11 + 1n0(1) + \partial(1)n1 + 2 \cdot (1(n + 1)1)\right)$$

$$= \partial(1(n + 1) + (n + 1)1 + 2 \cdot (1n1)) - (21n + 12n + 2 \cdot (111n) - 2n1 - 1(n + 1)1 - 2 \cdot (11n1).$$
i.e. \( \langle u, v \rangle - \langle v, u \rangle \) where in the second step we expanded the derivation terms \( \mathcal{D}(1n + n1) \), and for \( u, v \in \mathbb{R}(\mathcal{C})_* \), the expression \([u, v]_* \) denotes the Pöppe product commutator of \( u \) and \( v \), i.e. \([u, v]_* := u * v - v * u \). Using Lemma 6 with \( u \) and \( v \) as the empty compositions and \( n = 1 \) and \( b = n \), as well as vice-versa, we deduce that,

\[
1 * n - n * 1 = \mathcal{D}(1n - n1) \quad \Leftrightarrow \quad 1n - n1 = \mathcal{D}^{-1}(1 * n - n * 1).
\]

Substituting this into the expression involving \(-[1, 1n]_* + [1, n1]_* \equiv -[1, 1n - n1]_*\) just above and rearranging, we observe,

\[
1(n + 1) + (n + 1)1 + 2 : (1n1) = \mathcal{D}^{-1}\left(1 * \mathcal{D}(n) + \mathcal{D}(n) * 1\right) + \mathcal{D}^{-1}[1, \mathcal{D}^{-1}[1, n]_*]_*.
\]

By substituting this result into the expression for \((n + 2) \in \mathbb{R}(\mathcal{C})_*\) just preceding Lemma 6 we have established the first part of the following result. The second part is established in the proof of the theorem that follows.

**Theorem 5 (Non-Commutative Lax hierarchy)** Let \( n \in \mathbb{N} \) be a given odd natural number with \( n \geq 3 \). Assume there exists a polynomial \( \tilde{\pi}_n = \tilde{\pi}_n(1, 2, \ldots, (n - 2)) \) such that for the composition element \( n \in \mathbb{R}(\mathcal{C})_* \) we have

\[
n = n + \tilde{\pi}_n(1, 2, \ldots, (n - 2)).
\]

For \( n = 3 \), such a polynomial \( \tilde{\pi}_3 = -3 \cdot 1 * 1 \) exists. By successively applying the derivation operation twice we deduce

\[
(n + 2) = (n + 2) + \mathcal{D}^2\tilde{\pi}_n - (1 * n + n * 1) - \mathcal{D}^{-1}(1 * \mathcal{D}(n) + \mathcal{D}(n) * 1) - \mathcal{D}^{-1}[1, \mathcal{D}^{-1}[1, n]_*]_*.
\]

In addition there exists a polynomial \( \bar{\pi}_n = \bar{\pi}_n(1, 2, \ldots, n) \) in \( \mathbb{R}(\mathcal{C})_* \) such that

\[
\mathcal{D}^{-1}(1 * \mathcal{D}(n) + \mathcal{D}(n) * 1) + \mathcal{D}^{-1}[1, \mathcal{D}^{-1}[1, n]_*]_* = \bar{\pi}_n(1, 2, \ldots, n).
\]

**Remark 12** Note we can substitute for \( 'n' \) in the term \((1 * n + n * 1)\) in the second statement in Theorem 5 from the first statement in the theorem.

**Remark 13** The expression for \((n + 2)\) in Theorem 5 matches that for the non-commutative potential Korteweg-de Vries hierarchy in Carillo and Schoenlieb [10] Sec. II.

**Proof** We establish the existence of the polynomial \( \bar{\pi}_n = \bar{\pi}_n(1, 2, \ldots, n) \) for the terms shown in the theorem. From the arguments preceding the theorem, we know that

\[
\mathcal{D}^{-1}(1 * \mathcal{D}(n) + \mathcal{D}(n) * 1) + \mathcal{D}^{-1}[1, \mathcal{D}^{-1}[1, n]_*]_* \equiv 1(n + 1) + (n + 1)1 + 2 : (1n1). \]

Hence, equivalently, we can show there exists a polynomial \( \bar{\pi}_n = \bar{\pi}_n(1, 2, \ldots, n) \) such that

\[
1(n + 1) + (n + 1)1 + 2 : (1n1) = \bar{\pi}_n(1, 2, \ldots, n).
\]
The argument required to establish this is analogous to that used in the proof of Theorem 4 in Section 6 with some slight modifications, particularly in the first steps. Suppose \( \hat{\pi}_n = \hat{\pi}_n(1, 2, \ldots, n) \) has the form

\[
\hat{\pi}_n := \frac{1}{2(n+3)} \sum_{k=1}^{r(n+3)} \sum_{a_1a_2 \cdots a_k \in \mathcal{C}^*(n+2)} c_{a_1a_2 \cdots a_k} \cdot a_1 \ast a_2 \ast \cdots \ast a_k.
\]

Note, implicitly, we exclude the possibility of the single term monomial \('n + 2'\).

When we explicitly evaluate the Pöppe products in the monomials shown on the right in this expression for \( \hat{\pi}_n \), we generate linear combinations of compositions in \( \mathcal{C}(n + 2) \). Just like in the proof of Theorem 4, imagine constructing a table whose rows are equivalent to those shown in Table 2, once we neglect the final non-homogenous linear equation \( \hat{\pi}_n \). The remaining information it carries. The remaining block whose rows correspond to the compositions \( n2, (n+1)3, \ldots, 1(n+1) \), contains \( (n-1) \) homogeneous linear equations for \( \ell = 2, \ldots, n \), given by

\[
\hat{c}_{(n-\ell+2)(\ell-1)} + \hat{c}_{(n-\ell+1)\ell} = 0,
\]

and a final non-homogenous linear equation \( \hat{c}_{1n} = 1 \). Note, all the signature coefficients are equal to 1. Due to the inherent symmetry in the relation between the terms given by \( 1(n+1) + (n+1)1 + 2 \cdot (1n1) \) and \( \hat{\pi}_n \) proposed above, we replace this final equation by the homogenous linear equation \( -\hat{c}_{1n} + \hat{c}_{1n} = 0 \). The homogeneous linear system of equations represent by those shown above for \( \ell = 2, \ldots, n \) can be straightforwardly solved, systematically working through the equations in order, giving \( \hat{c}_{(n-\ell+1)\ell} = \).
The entries in the corresponding columns in the rows below this block to be zero.

We henceforth retain the homogeneous linear system of equations for \( \ell = 2, \ldots, n \) shown above, and the final replacement homogeneous linear equation. We nominate the “free” variable to be \( \tilde{c}_{n1} \). This resulting \( n \times n \) system of homogeneous linear equations is upper triangular with all the leading diagonal entries non-zero. Thus all these leading diagonal entries can be used as pivots in a Gaussian elimination process to render all the entries in the corresponding columns in the rows below this block to be zero.

We now adopt a strategy for solution for the linear system \( \tilde{A} \tilde{C} = \tilde{B} \), analogous to that in the proof of Theorem 4. However, before proceeding there is one minor snag in the procedure we must deal with first. The snag is that all the remaining linear equations are homogeneous apart from one, namely that for the row corresponding to the composition \( 1n1 \). This non-homogeneous linear equation has the form

\[
\chi_{(n-1)\otimes 1} \tilde{c}_{n1} + \chi_{1\otimes (n-1)1} \tilde{c}_{1n} + \chi_{1\otimes (n-2)\otimes 1} \tilde{c}_{1(n-2)1} = 2.
\]

Note hereafter in this proof we denote the arguments of the signature character maps \( \chi \) as sub-indices, as indicated in the equation above, for brevity. Since we know that \( \chi_{(n-1)\otimes 1} = \chi_{1\otimes (n-1)1} = \chi_{1(n-1)} = n \) and \( \chi_{1\otimes (n-2)\otimes 1} = 1 \), the equation above is equivalent to \( n \tilde{c}_{n1} + n \tilde{c}_{1n} + \tilde{c}_{1(n-2)1} = 2 \). This equation is in turn equivalent to the relation \( \tilde{c}_{1(n-2)1} = 2 - 2n \tilde{c}_{1n} \). We now demonstrate that this equation is consistent with the preceding pivot equations, in terms of descent order. Recall our dicussion involving the level 2, 3-part composition sub-blocks of \( A \) in the proof of Theorem 4. The pivot equation in the row preceding the row \( 1n1 \) for the equation shown above occurs in the row corresponding to the composition \( 2(n-1)1 \). The pivot equation in that row, i.e. in the row corresponding to \( 2(n-1)1 \), involves the unknown \( \tilde{c}_{2(2-3)1} \), and so we must invoke the pivot equation corresponding to that variable that lies in a preceding row. Iteratively continuing this procedure we arrive at the following system of homogeneous linear equations respectively given by the rows corresponding to the compositions \( 2(n-1)1, 3(n-2)1, 4(n-3)1, \ldots, (n-2)11 \):

\[
\begin{align*}
\tilde{c}_{1(n-2)1} &+ \tilde{c}_{2(n-3)1} + \chi_{(n-2)\otimes 1} \tilde{c}_{n1} + \chi_{2\otimes (n-2)1} \tilde{c}_{2(n-1)} + \chi_{1\otimes (n-1)1} \tilde{c}_{1n} = 0, \\
\tilde{c}_{2(n-3)1} &+ \tilde{c}_{3(n-4)1} + \chi_{(n-3)\otimes 1} \tilde{c}_{n1} + \chi_{3\otimes (n-3)1} \tilde{c}_{3(n-2)} + \chi_{2\otimes (n-2)1} \tilde{c}_{2(n-1)} = 0, \\
\tilde{c}_{3(n-4)1} &+ \tilde{c}_{4(n-5)1} + \chi_{(n-4)\otimes 1} \tilde{c}_{n1} + \chi_{4\otimes (n-4)1} \tilde{c}_{4(n-3)} + \chi_{3\otimes (n-3)1} \tilde{c}_{3(n-2)} = 0, \\
&\vdots
\end{align*}
\]

We now utilise that for \( \ell = 2, \ldots, n \) we know \( \tilde{c}_{(n-\ell+1)\ell} = (-1)^{(\ell-1)} \tilde{c}_{n1} \). Solving the homogeneous linear system of equations above in reverse we observe:

\[
\begin{align*}
\tilde{c}_{(n-2)\otimes 1} &= (\chi{(n-1)\otimes 1} - 1) \tilde{c}_{n1}, \\
\tilde{c}_{(n-3)\otimes 21} &= (\chi{(n-1)\otimes 1} - \chi{(n-2)\otimes 2} + 2) \tilde{c}_{n1}, \\
&\vdots \\
\tilde{c}_{1(n-2)\otimes 1} &= (\chi{(n-1)\otimes 1} + \cdots + (-1)^{(n-1)} \chi{(n-2)\otimes 2} + (-1)^{(n-2)} (n - 2)) \tilde{c}_{n1}.
\end{align*}
\]
We observe that the coefficient of \( \tilde{c}_{n1} \) on the right for odd \( n \) is given by,
\[
\chi_{n-1} - \chi_{n-2} + \cdots + (-1)^{n-1} \chi_{2(n-2)} + (-1)^{n-2}(n-2) \\
= \left( \frac{n}{n-1} \right) - \left( \frac{n}{n-2} \right) + \cdots + (-1)^{n-1} \left( \frac{n}{2} \right) + (-1)^{n-2}(n-2) \\
= - \left[ (1-x)^n - 1 - n(-x)^{n-1} - (-x)^n \right] \mid_{x=1} - (n-2) \\
= 2 - 2n.
\]
Hence we observe that we have the two expressions resulting from the row corresponding to the composition \( 1n1 \), namely \( \tilde{c}_{1(n-2)1} = 2 - 2n \tilde{c}_{1n} \) and \( \tilde{c}_{1(n-2)1} = (2 - 2n) \tilde{c}_{1n} \). Temporarily holding these two equations back, we observe that we can proceed to determine all of the coefficients \( \tilde{c}_{w} \) for \( w \in \mathbb{C}^* (n+1) \setminus (n+2) \) in terms of \( \tilde{c}_{n1} \) by precisely following the procedure outlined in the proof of Theorem 4. Finally we recover the first equation, i.e. that \( \tilde{c}_{n1} = 1 \), which determines all the coefficients \( \tilde{c}_{w} \) uniquely, and shows the two equations we temporarily held back are consistent.

Remark 14 (The Lax hierarchy is unique) The hierarchy given in Theorem 4 at each odd order \( n \) has a field given by a unique polynomial \( \pi_n = \pi_n(1, 2, \ldots, (n-2), n) \) in \( \mathbb{R} \langle \mathbb{C} \rangle \). For the Lax hierarchy presented in Theorem 5 the vector field has the polynomial form \( n+\hat{\pi}_n(1, 2, \ldots, (n-2)) \) in \( \mathbb{R} \langle \mathbb{C} \rangle \). The uniqueness of the former polynomial guarantees the two hierarchies are uniquely one and the same.

8 Conclusion

We have shown that integrability in the sense of linearisation for the whole Korteweg–de Vries hierarchy is equivalent to the existence of polynomial expansions for basic compositions in the real algebra of compositions equipped with the Pöppe product, \( \mathbb{R} \langle \mathbb{C} \rangle \). This procedure opens many new research directions we intend to pursue next:

(i) Do all compositions and/or linear combinations of compositions in \( \mathbb{R} \langle \mathbb{C} \rangle \) have Pöppe polynomial expansions?

(ii) There are natural connections to shuffle algebras, Rota–Baxter operators, symmetric functions and so forth, do these provide any further insight into integrability in this context?

(iii) The Pöppe product can be considered as a quasi-Leibniz product with the third term in Definition 7 considered to be the modification from the Leibniz product. This term is analogous to the the modification distinguishing the quasi-shuffle product from the shuffle product, see for example Curry et al. [17]. A natural question is, how robust are the results established in Sections 6 and 7 to perturbation of the coefficient ‘2’ in front of the third modification term in the Pöppe product?

(iv) The shuffle product and de-shuffle co-product are adjoint operators, as are the concatenation product and deconcatenation, see Reutenauer [60, p. 27]. Can we prove that the Pöppe product and de-Pöppe co-product are adjoint operations?

(v) Following on from (iv), and as hinted at in Section 5, a natural progression is to introduce the following tensor product of algebras: \( S \otimes \mathbb{R} \langle \mathbb{C} \rangle \). Note we have \( S \cong \mathbb{R} \langle \mathbb{C} \rangle \). The idea is to pull back linear combinations and/or polynomials in \( \mathbb{R} \langle \mathbb{C} \rangle \) to \( S \otimes \mathbb{R} \langle \mathbb{C} \rangle \) and then perform computations in this latter algebra. Then further natural questions arise concerning whether it is possible to
endow $S$ with a compatible product and $\mathbb{R}(C^\infty)$ a compatible co-product to raise their status to bi-algebras and perhaps even Hopf algebras?

(vi) The generation of the Korteweg–de Vries hierarchy via the Lax iteration shown in Theorem 5 involves evaluating a second order derivation $\partial^2$ together with the other terms shown. We know from Theorem 5 that overall, the terms generated via the Lax iteration can be expressed in the form of a polynomial in the signature expansions. Further for example, the nonlinear terms can be expressed in terms of derivations as follows. For the quintic non-commutative potential Korteweg–de Vries equation, the quadratic nonlinear terms can be encoded as $5 \cdot (\partial^2(1+1) - 2 \cdot 2)$. For the septic version in the hierarchy the quadratic nonlinear terms can be encoded as $7 \cdot (\partial^4(1+1) - 2 \cdot \partial^2(2+2) + 3 \cdot 3)$. The question is, can the cubic and higher degree terms be similarly encoded and this formulation of the nonlinear terms be propagated to all orders?

(vii) Finally, crucially, Doikou et al. [21] prove that the non-commutative nonlinear Schrödinger and modified Korteweg–de Vries flows are Grassmannian flows and integrable in this sense. Indeed we show that nonlocal reverse time and reverse space-time versions of these equations are also Grassmannian flows. See Ablowitz and Musslimani [2], Gürses and Pekcan [39,40,41,42] and Fokas ?? for more details on such systems. Further, Malham [46] demonstrated that the local and nonlocal non-commutative fourth order quintic nonlinear Schrödinger flow is also a Grassmannian flow. A natural extension is to determine if the combinatorial algebraic procedure we develop herein can be extended, in modified form, to all orders in the classical hierarchies for such non-commutative integrable systems? A modification is necessary as the linear system of equations required to prescribe the Grassmannian flow for these systems has the form: $\mathcal{P} = \mathcal{G}(\text{id} + Q)$, where

$$\mathcal{P} = \begin{pmatrix} P \\ \tilde{P} \end{pmatrix}, \quad Q = \begin{pmatrix} Q & O \\ O & Q \end{pmatrix} \quad \text{and} \quad \mathcal{G} = \begin{pmatrix} G \\ \tilde{G} \end{pmatrix}.$$

Here the block operators shown satisfy the following linear system of equations: $\partial_t P = d(\partial) P$, $\partial_t \tilde{P} = \tilde{d}(\partial) \tilde{P}$, $Q := \tilde{P} P$ and $\tilde{Q} := P \tilde{P}$, and $d$ as well as $\tilde{d}$ are suitable constant coefficient polynomials in $\partial$. Note once the operators $P$ and $\tilde{P}$ are determined as solutions to the linear partial differential systems shown, the operators $Q$ and $\tilde{Q}$ are prescribed by the quadratic forms shown. It is these definitions for $Q$ and $\tilde{Q}$ that necessitate, for the fully general case, a non-trivial extension of the combinatorial algebra introduced herein.

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We prove Theorem 3. First, let us prove axiom (i). Since the co-unit given in Definition 11 sends all non-empty words to 0 if we apply id \( \otimes \varepsilon \) to the form for \( \Delta(a_1 a_2 \cdots a_n) \) given in Lemma 13 we observe the only non-zero image would be from the final term for which (id \( \otimes \varepsilon \) \( \circ \) \( \varepsilon \otimes a_1 \cdots a_n \)) = id \( \circ \) \( a_1 \cdots a_n \). Likewise the only non-zero image of \( \varepsilon \otimes \text{id} \) would be from the first term for which \( (\varepsilon \otimes \text{id}) \circ (\nu \otimes a_1 \cdots a_n) = a_1 \cdots a_n \) = id \( \circ \) \( a_1 \cdots a_n \). Hence axiom (i) is satisfied.

Second, we focus on proving axiom (ii). It helps to consider the co-product \( \delta: \mathcal{C} \rightarrow \mathcal{C} \otimes \mathcal{C} \) which denotes the standard deconcatenation co-product, i.e. for any composition \( a_1 a_2 \cdots a_n \in \mathcal{C} \) we have

\[
\delta(a_1 a_2 \cdots a_n) = \sum_{k=0}^{n} a_1 \cdots a_k \otimes a_{k+1} \cdots a_n,
\]

where we use the convention that when \( k = 0 \) and \( k = n \), the corresponding terms are respectively \( \nu \otimes a_1 a_2 \cdots a_n \) and \( a_1 \cdots a_n \otimes \nu \). Note we have \( \delta(\nu) = \nu \otimes \nu \). We observe that if we apply (id \( \otimes \delta \)) to the last expression we get

\[
(id \otimes \delta) \circ \delta(a_1 a_2 \cdots a_n) = \sum_{k=0}^{n} a_1 \cdots a_k \otimes \delta(a_{k+1} \cdots a_n)
= \sum_{k=0}^{n} a_1 \cdots a_k \otimes \left( \sum_{\ell=0}^{n-k} a_{k+1} \cdots a_{k+\ell} \otimes a_{k+\ell+1} \cdots a_n \right)
= \sum_{k=0}^{n} \sum_{\ell=0}^{n-k} a_1 \cdots a_k \otimes a_{k+1} \cdots a_{k+\ell} \otimes a_{k+\ell+1} \cdots a_n.
\]

We observe the last line represents the sum over all possible ways to deconcatenate \( a_1 \cdots a_n \) into three parts, by performing a “left-to-right” secondary deconcatenation after the first. Hence we expect the action of \( (\delta \otimes \text{id}) \) on \( \delta(a_1 a_2 \cdots a_n) \) to correspond to performing “right-to-left” secondary deconcatenation after the first, producing the same triple deconcatenation result. Indeed we observe,

\[
(\delta \otimes \text{id}) \circ \delta(a_1 a_2 \cdots a_n) = \sum_{k=0}^{n} \delta(a_1 \cdots a_k) \otimes a_{k+1} \cdots a_n
= \sum_{k=0}^{n} \left( \sum_{\ell=0}^{k} a_1 \cdots a_\ell \otimes a_{\ell+1} \cdots a_k \right) \otimes a_{k+1} \cdots a_n
= \sum_{k=0}^{n} \sum_{\ell=0}^{k} a_1 \cdots a_\ell \otimes a_{\ell+1} \cdots a_k \otimes a_{k+1} \cdots a_n,
\]
\[
\begin{align*}
\sum_{\ell=0}^{n} \sum_{k=\ell}^{n} a_1 \cdots a_{\ell} \otimes a_{\ell+1} \cdots a_k \otimes a_{k+1} \cdots a_n &= \\
\sum_{\ell=0}^{n} \sum_{k=0}^{n-\ell} a_1 \cdots a_{\ell} \otimes a_{\ell+1} \cdots a_{k+\ell} \otimes a_{k+\ell+1} \cdots a_n,
\end{align*}
\]

which is the same as the “left-to-right” sum once we swap the identities of the dummy summation variables \(k\) and \(\ell\). The identification of \((\text{id} \otimes \Delta) \circ \Delta\) and \((\Delta \otimes \text{id}) \circ \Delta\) follows in exactly the same way. For an \(\Delta(a_1 \cdots a_n)\) co-product expansion, except for the \(\nu \otimes a_1 \cdots a_n\) and \(a_1 \cdots a_n \otimes \nu\) terms, we replace the terms \(a_1 \cdots a_k \otimes a_{k+1} \cdots a_n\) in \(\delta(a_1 \cdots a_n)\) for \(k = 1, \ldots, n-1\) by the pair of terms

\[
a_1 \cdots a_{k-1} (a_k - 1) \otimes a_{k+1} \cdots a_n + a_1 \cdots a_k \otimes (a_{k+1} - 1) a_{k+2} \cdots a_n.
\]

The co-product \(\Delta(a_1 \cdots a_n)\) tensorally decomposes \(a_1 \cdots a_n\) into the sum of all possible composition pairs that produce \(a_1 \cdots a_n\) via the Pöpppe product \(\ast\), including the empty composition \(\nu\). Hence \((\text{id} \otimes \Delta) \circ \Delta\) will tensorally decompose \(a_1 \cdots a_n\) into the sum of all possible composition triples \(u \otimes v \otimes w\) that produce \(a_1 \cdots a_n\) via the Pöpppe product \(u \ast v \ast w\). The action of \((\text{id} \otimes \Delta) \circ \Delta\) is to perform a “left-to-right” secondary de-Pöpppe co-product after a first application of the de-Pöpppe co-product. The action of \((\Delta \otimes \text{id}) \circ \Delta\) is to perform a “right-to-left” secondary de-Pöpppe co-product after a first application of the de-Pöpppe co-product. Both generate the same sum over all triples \(u \otimes v \otimes w\) that produce \(a_1 \cdots a_n\). As in the proof of Lemma 5, care must be taken to account for instances when \(a_k = 1\). These can be taken care of by a post-application of \(\theta \otimes \theta \otimes \theta\) and extending the definition of \(\theta\) to include the circumstance \(\theta: -1 \mapsto 0 \nu\) where in the image \(0 \in R\) so the corresponding term is eliminated. The expression \(-1\) could arise through a double application of \(\delta^{-1}\) to a letter \(a = 1\). We have thus established that axiom (ii) is satisfied. The proof is complete.
Table 3 Non-zero signature coefficients in $\pi_n$ appearing in the 3-part composition blocks of $A$. The coefficients are the $\chi$-images of the signature entries shown.

| $C$     | $(n-4)$-1-1 | $(n-5)$-2-1 | $(n-5)$-1-2 | $(n-6)$-3-1 | $(n-6)$-2-2 | $(n-6)$-1-3 | $(n-7)$-4-1 | $(n-7)$-3-2 | $(n-7)$-2-3 | $(n-7)$-1-4 | $(n-8)$-5-1 |
|---------|-------------|-------------|-------------|-------------|-------------|-------------|-------------|-------------|-------------|-------------|-------------|
| $(n-2)11$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-3)21$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-3)12$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-4)31$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-4)22$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-4)13$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-5)41$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-5)32$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-5)23$ | $(n-5)\otimes 1 \otimes 1$ | $(n-5)\otimes 1 \otimes 1$ | $(n-5)\otimes 1 \otimes 1$ | $(n-5)\otimes 1 \otimes 1$ | $(n-5)\otimes 1 \otimes 1$ | $(n-5)\otimes 1 \otimes 1$ | $(n-5)\otimes 1 \otimes 1$ | $(n-5)\otimes 1 \otimes 1$ | $(n-5)\otimes 1 \otimes 1$ | $(n-5)\otimes 1 \otimes 1$ | $(n-5)\otimes 1 \otimes 1$ |
| $(n-6)51$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-6)42$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-6)33$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-6)24$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-7)61$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-7)52$ |             |             |             |             |             |             |             |             |             |             |             |
| $(n-7)43$ |             |             |             |             |             |             |             |             |             |             |             |
| ...      |             |             |             |             |             |             |             |             |             |             |             |
Table 4: Non-zero signature coefficients in \( \tau_n \) appearing in the 4-part composition blocks of \( A \). The coefficients are the \( \chi \)-images of the signature entries shown.