Modified poisson series to implement the method of multiple scales

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The method of multiple scales is a global perturbation technique that has resulted to be very useful in perturbed ordinary differential equations characterized by disparate time scales. The general principle behind the method is that the solution to the differential equation is uniformly expanded in terms of two or more independent variables, referred to as time scales. In this article, we present a mathematical object based on a Poisson series to apply the method of multiple scales via specific symbolic computation. Copyright © 2016 John Wiley & Sons, Ltd.

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1. Introduction

Perturbation theories for differential equations containing a small parameter $\epsilon$ are quite old. The small perturbation theory originated by Newton has been highly developed by many others, and an extension of this theory to the asymptotic expansion, consisting on a power series expansion in the small parameter, was devised by Poincaré [1]. The main point is that for the most of the differential equations, it is not possible to obtain an exact solution. In cases where equations contain a small parameter, we can consider it as a perturbation parameter to obtain an asymptotic expansion of the solution. These classical perturbation methods based on asymptotic expansions generally break down because of resonances that lead to what are called secular terms. Multiple scales analysis [2] comprises techniques used to construct uniformly valid approximations to the solutions of perturbation problems, both for small as well as large values of the independent variables. These techniques are useful in systems characterized by disparate time scales, such as weak dissipation in an oscillator. In the standard approach, the expansion of the solution to the perturbed equation depends explicitly on $t$, $\epsilon t$, $\epsilon^2 t$, and the small parameter $\epsilon$ itself. The method of multiple scales has become very popular, and it has been applied to a wide range of problems. To cite some examples, Cole and Kevorkian [3], Nayfeh [4–6], Musa [7] and Reiss [8] applied the method of multiple scales to the analysis of weakly linear and nonlinear vibrations modeled by second and third order ordinary differential equations. Rammath and Sandri [9] used it to study equations with variable coefficients. More recently, Krämer [10] has computed an approximate solution to the nonlinear Klein–Gordon equation via the method of multiple scales. Abbasi et al. [11] employ the method of multiple scales to analyze the chaotic behavior and different types of fixed points in ferroresonance of voltage transformers considering core loss. This phenomenon has nonlinear chaotic dynamics and includes subharmonic, quasi-periodic, and also chaotic oscillations. In [12] and [13], the authors generalise a computer implementation of the multiple scales method and its application to nonlinear vibration problems. The necessary macro-steps that are used for the development of the computational system are formulated, and the practical ways of encoding these steps using Mathematica are discussed. Cartmell et al. [14] investigate the application of the method to some problems in the area of machine and structural dynamics.

In practice, the work involved in the application of this approach to compute the solution to a differential equation cannot be performed by hand, and algebraic processors result to be a very useful tool. Since the early sixties, many systems for specific or general symbolic computation have been developed. Nowadays, many general purpose computer algebra packages contain tools for the calculation of the solution of certain classes of ordinary differential equations. All these packages have the advantage of being very general, so they can deal with a lot of problems of different nature. However, the most common perturbation methods tend to produce expressions containing thousands of terms, and their treatment with those general processors becomes a time-consuming task. Specific symbolic computation packages avoid this inconvenience working with simple data structures and algorithms.
2. The method of multiple scales

We will consider the initial value problem defined by the following nonlinear second order differential equation:

$$\ddot{x} + x = \epsilon f(x, \dot{x}),$$

where $0 < \epsilon \ll 1$ is a small parameter and $f(x, \dot{x})$ can be arranged as follows,

$$f(x, \dot{x}) = \sum_{q=0}^{M} \sum_{0 \leq v \leq q} f_{v,q} x^v \dot{x}^{q-v},$$

being $f_{v,q} \in \mathbb{R}$ for $0 \leq v, q \leq M$, and $M \in \mathbb{N}$. The standard approach [15] is to try a power series solution of the form

$$x(t, \epsilon) = x_0(t) + \epsilon x_1(t) + \epsilon^2 x_2(t) + \cdots.$$ \hspace{1cm} (3)

This series is inserted into the governing equation and initial conditions, and coefficients of same powers of $\epsilon$ are then grouped to obtain a collection of equations for the coefficient functions $x_i(t)$, which are then solved in a sequential manner. The resulting solution $x(t, \epsilon)$ can be useful in approximating the function when $\epsilon$ is small, but it results to be a poor approximation when $t$ is as large as $\epsilon^{-1}$. To determine an expansion valid for times as large as $\epsilon^{-1}$, the combination $\epsilon t$ should be considered a single variable $T_1 = \epsilon t$. Thus, the truncated expansion is valid for times as large as $\epsilon^{-1}$, but it is not satisfactory when $t$ is as large as $\epsilon^{-2}$. To obtain an adequate asymptotic expansion valid for $t = O(\epsilon^{-2})$, $\epsilon^2 t$ should be considered a single variable $T_2 = \epsilon^2 t$.

The aforementioned discussion suggests that $x(t, \epsilon)$ depends explicitly on $T_0 = t$, $T_1 = \epsilon t$, $T_2 = \epsilon^2 t$, and so on. In order to get a truncated expansion valid for any $t$ up to $O(\epsilon^n)$, we must calculate the dependence of $x$ on the $n+1$ different time scales $T_0, T_1, \ldots, T_n$. The method of multiple scales considers the expansion to be a function of multiple independent variables, or scales, instead of a single variable $t$,

$$x(t) = \sum_{v=0}^{n-1} \epsilon^v x_v(T_0, T_1, \ldots, T_n) + O(\epsilon^n),$$

where the independent variables are defined as

$$T_v = \epsilon^v t, \quad v = 0, 1, 2, \ldots, n.$$ \hspace{1cm} (4)

The number of independent time scales depends on the order to which the expansion is carried out. Substitution of Equation (3) into the governing differential equation and collecting coefficients of equal powers of $\epsilon$ generate a system of $n+1$ differential equations. To obtain a uniform solution, the system of differential equations must be solved sequentially for $v = 0, 1, \ldots, n-1$, eliminating secular terms.

The time derivative is transformed according to

$$\frac{d}{dt} = \frac{\partial}{\partial T_0} + \epsilon \frac{\partial}{\partial T_1} + \epsilon^2 \frac{\partial}{\partial T_2} + \cdots.$$ \hspace{1cm} (5)

For our purpose this time, we will focus our attention in the functions $x_i$ depending only on two time scales, $T_0 = t$ and $T_1 = \epsilon t$. Of course we can extend this procedure to as many time scales as we like, but because some upcoming equations would become unnecessarily difficult without any further insight into the method itself, we leave it at two time scales. Thus, considering two time scales, we get

$$x(t, \epsilon) = x_0(T_0, T_1) + \epsilon x_1(T_0, T_1) + O(\epsilon^2).$$

Taking into account the transformation of the time derivative,

$$\frac{d}{dt} = \frac{\partial}{\partial T_0} + \epsilon \frac{\partial}{\partial T_1},$$

we get

$$\dot{x} = \frac{\partial}{\partial T_0} x_0 + \epsilon \frac{\partial}{\partial T_1} x_0 + \epsilon \frac{\partial}{\partial T_0} x_1 + O(\epsilon^2).$$

Let us define the operators

$$D_0 = \frac{\partial}{\partial T_0}, \quad D_1 = \frac{\partial}{\partial T_1}.$$
Then, equations for \( \dot{x} \) and \( \ddot{x} \) can be written as

\[
\dot{x} = D_0 x_0 + \epsilon (D_1 x_0 + D_0 x_1) + O(\epsilon^2),
\]

\[
\ddot{x} = D_0 \ddot{x}_0 + \epsilon (2D_1 \dot{x}_0 + D_0 \ddot{x}_1) + O(\epsilon^2).
\]

If we substitute these equations into the governing Equation (1), we get

\[
D_0 \ddot{x}_0 + \epsilon (2D_0 \dot{x}_0 + D_0 \ddot{x}_1) + x_0 + \epsilon x_1 = \epsilon \sum_{q=0}^{M} \sum_{0 \leq \nu \leq q} x_0^{\nu} (D_0 x_0)^{q-\nu}.
\]

Now, we collect terms in equal powers of \( \epsilon \), to generate the following system of equations,

\[
O(1): D_0 \ddot{x}_0 + x_0 = 0,
\]

\[
O(\epsilon): D_0 \ddot{x}_1 + x_1 = \sum_{q=0}^{M} \sum_{0 \leq \nu \leq q} x_0^{\nu} (D_0 x_0)^{q-\nu} - 2D_0 \dot{x}_0.
\]

The solution to (1) is constructed from the order zero, which corresponds with the unperturbed problem, and can be written as

\[
D_0 \ddot{x}_0 + x_0 = 0.
\]

The solution to Equation (8) is

\[
x_0(T_0, T_1) = R(T_1) \cos (T_0 + \Phi(T_1)) ,
\]

where \( R(T_1) \) and \( \Phi(T_1) \) are the slow-varying amplitude and phase of \( x_0 \), respectively. Now, we need to compute \( D_0 \dot{x}_0 \) and \( D_1 \ddot{x}_0 \) in order to obtain \( x_1 \). Derivating Equation (9), we get

\[
D_0 \dot{x}_0 = -R(T_1) \sin (T_0 + \Phi(T_1)) ,
\]

and

\[
D_1 \ddot{x}_0 = -D_1 R(T_1) \sin (T_0 + \Phi(T_1)) - R(T_1) D_1 \Phi(T_1) \cos (T_0 + \Phi(T_1)) .
\]

Now, we can substitute Equations (10) and (11) into Equation (7). First, we need to calculate the products

\[
(x_0)^\nu = R^\nu (T_1) \cos^\nu (T_0 + \Phi(T_1)),
\]

\[
(D_0 x_0)^{q-\nu} = (-1)^{q-\nu} R^{q-\nu} (T_1) \sin^{q-\nu} (T_0 + \Phi(T_1)),
\]

but also

\[
-2D_1 \dot{x}_0 = 2D_1 R(T_1) \sin (T_0 + \Phi(T_1)) + 2R(T_1) D_1 \Phi(T_1) \cos (T_0 + \Phi(T_1)).
\]

Thus, Equation (7) can be arranged as follows,

\[
D_0 \ddot{x}_1 + x_1 = \sum_{q=0}^{M} \sum_{0 \leq \nu \leq q} f_{\nu, q-\nu} (-1)^{q-\nu} R^{q-\nu} (T_1) \sin^{q-\nu} (T_0 + \Phi(T_1)) + + 2D_1 R(T_1) \sin (T_0 + \Phi(T_1)) + 2R(T_1) D_1 \Phi(T_1) \cos (T_0 + \Phi(T_1)).
\]

Now, we must collect terms of same frequencies in order to avoid the generation of secular terms, by making resonant terms vanish. To do that, we must expand the right-hand side of Equation (12) by means of the relations

\[
\sin(\alpha_1 + \alpha_2) = \sin \alpha_1 \cos \alpha_2 + \sin \alpha_2 \cos \alpha_1 ,
\]

\[
\cos(\alpha_1 + \alpha_2) = \cos \alpha_1 \cos \alpha_2 - \sin \alpha_1 \sin \alpha_2 .
\]

In the following section, we introduce an example to show clearly the type of expressions we have to handle in order to construct a symbolic computation system to apply the method of multiple scales.
3. Application of the method to the van der Pol equation

In this section, we apply the two timing method to show the mathematical object we have to handle with the specific symbolic computation system. To that purpose, let us consider the van der Pol equation,

\[ \ddot{x} + x = \epsilon (1 - x^2) \dot{x}. \] (13)

The van der Pol oscillator is a non-conservative oscillator with nonlinear damping. Energy is dissipated at high amplitudes and generated at low amplitudes. Thus, there exists oscillations around a state at which energy generation and dissipation balance. The state towards which the oscillations converge is known as a limit cycle. Van der Pol discovered that no matter the initial conditions, the oscillator converges to a limit cycle. For \( \epsilon \ll 1 \) and trajectories close to the origin, the amplitude of the oscillation grows very slowly, each oscillation with a different amplitude and period. This behavior gives rise to the concept of multiple timescale of oscillation. The van der Pol oscillator has become the cornerstone for studying systems with limit cycle oscillations in physics, biology, sociology, and even economics.

The \( O(1) \) and \( O(\epsilon) \) equations ((6) and (7), respectively) become

\[ D_0 D_0 x_0 + x_0 = 0, \] (14)

and

\[ D_0 D_0 x_1 + x_1 = -x_0^2 D_0 x_0 + D_0 x_0 - 2D_0 D_1 x_0. \] (15)

The solution to (14) is

\[ x_0(T_0, T_1) = R(T_1) \cos(T_0 + \Phi(T_1)), \]

where \( R(T_1) \) and \( \Phi(T_1) \) are the slow-varying amplitude and phase of the oscillation. The derivatives \( D_0 x_0 \) and \( D_1 D_0 x_0 \) can be written as

\[ D_0 x_0 = -R(T_1) \sin(T_0 + \Phi(T_1)) \]

and

\[ D_1 D_0 x_0 = -D_1 R(T_1) \sin(T_0 + \Phi(T_1)) - R(T_1) D_1 \Phi(T_1) \cos(T_0 + \Phi(T_1)). \]

Now, we substitute the solution to (14) as well as its derivatives into (15) to get

\[ D_0 D_0 x_1 + x_1 = R^2(T_1) \cos^2(T_0 + \Phi(T_1)) \sin(T_0 + \Phi(T_1)) - R(T_1) \sin(T_0 + \Phi(T_1)) + 2D_1 R(T_1) \sin(T_0 + \Phi(T_1)) + 2R(T_1) D_1 \Phi(T_1) \cos(T_0 + \Phi(T_1)). \] (16)

Taking into account the relation

\[ \cos^2(T_0 + \Phi(T_1)) \sin(T_0 + \Phi(T_1)) = \frac{1}{4} \sin(T_0 + \Phi(T_1)) + \frac{1}{4} \sin(3T_0 + 3\Phi(T_1)), \]

we can expand Equation (16) to obtain

\[ D_0 D_0 x_1 + x_1 = \frac{1}{4} R^2(T_1) \sin(T_0 + \Phi(T_1)) + \frac{1}{4} R^2(T_1) \sin(3T_0 + 3\Phi(T_1)) - R(T_1) \sin(T_0 + \Phi(T_1)) + 2D_1 R(T_1) \sin(T_0 + \Phi(T_1)) + 2R(T_1) D_1 \Phi(T_1) \cos(T_0 + \Phi(T_1)). \] (17)

Now, we can collect terms in the same frequencies in order to avoid the generation of secular terms,

\[ D_0 D_0 x_1 + x_1 = \left( \frac{1}{4} R^2(T_1) - R(T_1) + 2D_1 R(T_1) \right) \sin(T_0 + \Phi(T_1)) + 2R(T_1) D_1 \Phi(T_1) \cos(T_0 + \Phi(T_1)) + \frac{1}{4} R^2(T_1) \sin(3T_0 + 3\Phi(T_1)). \] (18)

Thus, resonant terms vanish if we make

\[ H_1(T_1) = \frac{1}{4} R^2(T_1) - R(T_1) + 2D_1 R(T_1) = 0, \]

\[ H_2(T_1) = 2R(T_1) D_1 \Phi(T_1) = 0. \]

So the conditions to avoid secular terms take the form of two differential equations,

\[ 2D_1 R(T_1) = R(T_1) - \frac{1}{4} R^2(T_1) - R(T_1) + 2D_1 R(T_1), \]

\[ D_1 \Phi(T_1) = 0. \]
The right-hand side of both equations is a polynomial in \( R \) and \( \Phi \), so we can compute the zeros of both sides. In this case, we get that 
\( R(T_1) = 0, \pm 2 \) and \( \Phi(T_1) \) is a linear function of \( T_1 \). A numerical exploration is enough to show that \( R(T_1) = 2 \) is a stable fixed point of the corresponding differential equation.

4. Mathematical object

From the analysis of Section 3, we see that a modification of a Poisson series is enough to handle the expressions involved in the application of the method of multiple scales. In this section, we will study the way the structure of a Poisson series must be adapted in order to implement the method of multiple scales as a symbolic algorithm.

A Poisson series is an expression of the form

\[
S = \sum_{i \in \mathbb{I}} \sum_{j \in \mathbb{K}} C_{ij} \xi_i^1 \xi_i^2 \cdots \xi_i^m \cos \bigg( \sum_{j=1}^n \phi_j + \cdots + \phi_n \bigg) \sin \bigg( \sum_{j=1}^n \phi_j + \cdots + \phi_n \bigg) .
\]  

(19)

The summations bear on collections \( \mathbb{I} \) and \( \mathbb{K} \) consisting of vector of integers \( i = (i_1, i_2, \ldots, i_m) \) and \( j = (j_1, j_2, \ldots, j_n) \), the coefficients \( C_{ij} \) belong to \( \mathbb{B} \). The term \( \xi_i^1 \xi_i^2 \cdots \xi_i^m \) is usually referred to as monomial, and the letters \( A_1, A_2, \ldots, A_m \) as the polynomial variables. The greek letters \( \Phi_1, \Phi_2, \ldots, \Phi_n \) constitute the angle variables.

The set of Poisson series forms a commutative algebra over the ring of coefficients \([16]\). If \( P \) and \( Q \) are Poisson series, then their sum and product by a real number, \( P + Q \) and \( aP \), with \( a \in \mathbb{R} \), are also Poisson series. Algebraic closure properties make automatic manipulation rather easy when the elements in the algebra are represented in a standard canonical form, because closure implies that the result retains the standard form of the operands.

Now, we will proceed to introduce some elements in the Poisson series structure at the computational level in order to be able to apply the method of multiple scales. Let \( \mathbb{T} \) be the set of independent variables (time scales) \( T_0, T_1, \ldots, T_n \) that is,

\[
\mathbb{T} = \{ T_0, T_1, \ldots, T_n \} ,
\]

and \( \mathbb{P} \) be the power set of \( \mathbb{T} \), \( \mathbb{P} = \mathcal{P}(\mathbb{T}) \). For example, if we consider \( n = 1 \), then \( \mathbb{T} = \{ T_0, T_1 \} \), with \( T_0 = t, T_1 = \epsilon t \), and \( \mathbb{P} = \{ \emptyset, \{ T_0 \}, \{ T_1 \}, \{ T_0, T_1 \} \} \). We will use these sets to establish how the integration constants depend on the different time scales.

Now, let us consider the \( m \) variables \( A_1(T_1), A_2(T_2), \ldots, A_m(T_m) \), where \( T_1, \ldots, T_m \in \mathbb{P} \). In the example earlier, \( n = 1, m = 2, A_1(T_1) = R(T_1) \) and \( A_2(T_2) = \Phi(T_1) \), that is, \( T_1 = T_2 = T \).

In order to apply the method of multiple scales as described in the example of Section 3, we also need to introduce the partial derivative of a variable with respect to a time scale in the symbolic object at the computational level. For that purpose, we first introduce the notation \( \sigma = (\sigma_2, \sigma_1) \), being \( \sigma_1, \sigma_2 \in \{-1, 0, 1, \ldots, n\} \). Then, we define the operator

\[
\partial_{\sigma_2, \sigma_1} = \frac{\partial^2}{\partial T_{\sigma_2} \partial T_{\sigma_1}}.
\]

If \( \sigma_1 \) or \( \sigma_2 \) equals \(-1\), the corresponding partial derivative is not performed.

In order to clarify how we introduce this operator, let us refer to the example described in Section 3, where \( n = 1, A_1(T_1) = R(T_1) \) and \( A_2(T_2) = \Phi(T_1) \). The solution to Equation \((14)\) is written as a symbolical object of the form

\[
x_0(T_0, T_1) = A_1(T_1) D_{-1,-1} A_1(T_1) \cos(T_0 + A_2(T_2)).
\]

The operator \( D_{-1,-1} \) implies that no partial derivative affects the variable \( A_1(T_1) \), so we do not have to take into account this part of the term, that is,

\[
D_{-1,-1} A_1(T_1) = 1.
\]

The derivative of \( x_0(T_0, T_1) \) with respect to \( T_0 \) is written as

\[
D_{0} x_0(T_0, T_1) = -A_1(T_1) D_{-1,-1} A_1(T_1) \sin(T_0 + A_2(T_2)).
\]

We see in the equation earlier that no partial derivative acts on \( A_1(T_1) \), as this variable does not depend on \( T_0 \). As before, \( D_{-1,-1} A_1(T_1) = 1 \). Now, the partial derivative of \( D_{0} x_0(T_0, T_1) \) with respect to \( T_1 \) reads

\[
D_{1} D_{0} x_0(T_0, T_1) = -D_{-1,1} A_1(T_1) \sin(T_0 + A_2(T_2)) - A_1(T_1) D_{-1,-1} A_1(T_1) D_{-1,1} A_2(T_2) \cos(T_0 + A_2(T_2)).
\]

Here, \( D_{-1,1} A_1(T_1) = 1 \), and

\[
D_{-1,1} A_1(T_1) = \frac{\partial}{\partial T_1} A_1(T_1), \quad D_{-1,1} A_2(T_2) = \frac{\partial}{\partial T_1} A_2(T_2).
\]
In order to handle the more general type of mathematical expression that appears in the application of the method of multiple scales as described in Sections 2 and 3, we will consider the following modified Poisson series,

\[ S = \sum_{\nu \in \mathbb{N}} P_{\nu} (A_1 (\tau_1), \ldots, A_m (\tau_m)) \sin (\nu_0 T_0 + \nu_1 A_1 (\tau_1) + \cdots + \nu_m A_m (\tau_m)) + \]

\[ + Q_{\nu} (A_1 (\tau_1), \ldots, A_m (\tau_m)) \cos (\nu_0 T_0 + \nu_1 A_1 (\tau_1) + \cdots + \nu_m A_m (\tau_m)), \]

where the summations bear on the collection \( \mathbb{N} \) consisting of vectors of integers \( \nu = (\nu_0, \nu_1, \ldots, \nu_m) \), and \( P_{\nu} \) and \( Q_{\nu} \) are expressions of the form

\[ P_{\nu} (A_1 (\tau_1), \ldots, A_m (\tau_m)) = \sum_{j \in \mathbb{K}} C^j_{\nu} A_1 (\tau_1) \times A_2 (\tau_2) \times \cdots A_m (\tau_m) D_{\sigma_{\nu}} A_1 (\tau_1) \times A_2 (\tau_2) \times \cdots A_m (\tau_m), \]

and

\[ Q_{\nu} (A_1 (\tau_1), \ldots, A_m (\tau_m)) = \sum_{j \in \mathbb{K}} C^j_{\nu} A_1 (\tau_1) \times A_2 (\tau_2) \times \cdots A_m (\tau_m) D_{\sigma_{\nu}} A_1 (\tau_1) \times A_2 (\tau_2) \times \cdots A_m (\tau_m). \]

Here, the summations bear on the collection \( \mathbb{K} \) consisting of vectors of integers \( \nu = (\nu_1, \ldots, \nu_m) \), and \( \sigma_{\nu_j} = (\sigma_{\nu_2}, \sigma_{\nu_3}) \), with \( \sigma_{\nu_2}, \sigma_{\nu_3} \in \{-1, 0, 1, \ldots, \nu_0\} \).

The letters \( A_1, A_2, \ldots, A_m \) are the polynomial and angle variables of the Poisson series. These variables depend on the time scales \( T_0, T_1, \ldots, T_n \) as given by \( \tau_1, \tau_2, \ldots, \tau_m \in \mathbb{R} \),

\[ A_1 = A_1 (\tau_1), \quad A_2 = A_2 (\tau_2), \quad \ldots \quad A_m = A_m (\tau_m). \]

The operators \( D_{\sigma_{\nu}} = D_{\sigma_{\nu_2}, \sigma_{\nu_3}} \), defined as

\[ D_{\sigma_{\nu_2}, \sigma_{\nu_3}} = \frac{\partial^2}{\partial T_{\sigma_{\nu_2}} \partial T_{\sigma_{\nu_3}}}, \]

for any \( \sigma_{\nu_1}, \sigma_{\nu_2} \neq -1 \), act on the corresponding variable \( A_{\nu} \), for any \( v = 1, \ldots, m \). These operators only act on the polynomial part of the term. If \( \sigma_{\nu_2} = -1 \) and \( \sigma_{\nu_3} \neq -1 \), then

\[ D_{\sigma_{\nu_2}, \sigma_{\nu_3}} (A_{\nu} (\tau_\nu)) = 1. \]

On the other hand, if \( \sigma_{\nu_1} = \sigma_{\nu_2} = -1 \), then

\[ D_{\sigma_{\nu_1}, \sigma_{\nu_2}} (A_{\nu} (\tau_\nu)) = 1. \]

Now, we look for a canonical representation for each equivalence class defined in the set of modified Poisson series [16]. For that purpose, the following operations must be performed over each series:

1. Let us consider a modified Poisson series

\[ S = \sum_{\nu \in \mathbb{N}} P_{\nu} (A_1 (\tau_1), \ldots, A_m (\tau_m)) \sin (\nu_0 T_0 + \nu_1 A_1 (\tau_1) + \cdots + \nu_m A_m (\tau_m)) + \]

\[ + Q_{\nu} (A_1 (\tau_1), \ldots, A_m (\tau_m)) \cos (\nu_0 T_0 + \nu_1 A_1 (\tau_1) + \cdots + \nu_m A_m (\tau_m)). \]

If \( \nu_0 < 0 \), the following rules must be applied:

\[ \sin (\nu_0 T_0 + \alpha) = -\sin (-\nu_0 T_0 - \alpha), \quad \cos (\nu_0 T_0 + \alpha) = \cos (-\nu_0 T_0 - \alpha). \]

Moreover, terms with identical trigonometric part must be grouped together.

2. The terms of a modified Poisson series must be ordered in the following way. Let us consider two terms of a modified Poisson series,

\[ S_{v} = P_{\nu} (A_1 (\tau_1), \ldots, A_m (\tau_m)) \sin (\nu_0 T_0 + \nu_1 A_1 (\tau_1) + \cdots + \nu_m A_m (\tau_m)) + \]

\[ + Q_{\nu} (A_1 (\tau_1), \ldots, A_m (\tau_m)) \cos (\nu_0 T_0 + \nu_1 A_1 (\tau_1) + \cdots + \nu_m A_m (\tau_m)), \]

and

\[ S_{v} = P_{\nu} (A_1 (\tau_1), \ldots, A_m (\tau_m)) \sin (\alpha_0 T_0 + \alpha_1 A_1 (\tau_1) + \cdots + \alpha_m A_m (\tau_m)) + \]

\[ + Q_{\nu} (A_1 (\tau_1), \ldots, A_m (\tau_m)) \cos (\alpha_0 T_0 + \alpha_1 A_1 (\tau_1) + \cdots + \alpha_m A_m (\tau_m)). \]

\[ \end{equation}
We say that $S_v < S_u$ if for the first integer $i \in \mathbb{Z}$, $0 \leq i \leq m$, such that $v_i \neq u_i$, then $v_i < u_i$.

3. The polynomials $P_v$ and $Q_v$, of each term of the modified Poisson series must be ordered as follows. Let

$$P_v(A_1(t_1), \ldots, A_m(t_m)) = \sum_{i \in \mathbb{R}} C_{\alpha_i} A_1^{\alpha_1}(t_1)D_{\sigma_i} A_1(t_1) \times A_2^{\alpha_2}(t_2)D_{\sigma_2} A_2(t_2) \times \cdots \times A_m^{\alpha_m}(t_m)D_{\sigma_m} A_m(t_m),$$

be a polynomial part of a term of a modified Poisson series. Here, $i = (i_1, \ldots, i_m) \in \mathbb{R}$, and $\sigma_i = (\sigma_{i_1}, \ldots, \sigma_{i_m})$, with $\sigma_{i_1}, \ldots, \sigma_{i_m} \in \{0, 1, \ldots, n\}$.
Let us consider now two terms of this polynomial part,

$$B_i = C_{\alpha_i} A_1^{\alpha_1}(t_1)D_{\sigma_i} A_1(t_1) A_2^{\alpha_2}(t_2)D_{\sigma_2} A_2(t_2) \times \cdots \times A_m^{\alpha_m}(t_m)D_{\sigma_m} A_m(t_m),$$

and

$$B_j = C_{\alpha_j} A_1^{\alpha_1}(t_1)D_{\sigma_j} A_1(t_1) A_2^{\alpha_2}(t_2)D_{\sigma_2} A_2(t_2) \times \cdots \times A_m^{\alpha_m}(t_m)D_{\sigma_m} A_m(t_m).$$

For the sake of simplicity, let us introduce the vector of integers

$$\zeta = (\zeta_1, \ldots, \zeta_m) \in \Omega,$$

such that

$$\zeta_1 = i_1, \quad \zeta_2 = i_2, \quad \zeta_3 = \sigma_{i_1}, \quad \zeta_4 = i_2, \quad \zeta_5 = \sigma_{i_2}, \quad \zeta_6 = \sigma_{i_3},$$

and so on. In general,

$$\zeta_{3p+1} = i_{p+1}, \quad \zeta_{3p+2} = \sigma_{i_{p+1}}, \quad \zeta_{3p+3} = \sigma_{i_{p+1}},$$

for any integer $p$ such that $0 \leq p < m$. Let $\Omega$ be the collection of vectors of integers of the form $\zeta = (\zeta_1, \ldots, \zeta_m)$. Let us consider now the corresponding vectors of indices $\zeta^{(i)}$ and $\zeta^{(j)}$, associated to terms $B_i$ and $B_j$, respectively. We say that $B_i < B_j$ if for the first integer $k$ such that $\zeta_k^{(i)} \neq \zeta_k^{(j)}$, then $\zeta_k^{(i)} < \zeta_k^{(j)}$ is satisfied.

5. Data structure for a modified Poisson series

Now, we will consider the special Poisson series set we are working with from the computational point of view. To that purpose, we will analyze the basic information that characterizes a modified Poisson series, as well as the data structure to store it in the computer. This must be done preserving the canonical representation we have chosen.

First of all, we must set the constants $m$ and $n$, characterizing the number of variables $A_1, \ldots, A_m$, and the number of time variables, $T_0, T_1, \ldots, T_m$, respectively. Once these constants are set, let us consider the $m$ variables $A_1(t_1), A_2(t_2), \ldots, A_m(t_m)$, where $t_1, \ldots, t_m \in \mathbb{R}$.

We need to represent the dependence of each variable $A_v$ on the set of independent variables $T_0, T_1, \ldots, T_m$ through $T_v$, for any $v = 1, \ldots, m$. To that purpose, we will use a $m \times (n + 1)$ matrix, $M_v$. Each row $v$ of the matrix indicates the dependence of the corresponding variable $A_v$ on the time variables. Thus,

$$(M_v)_{v, \alpha} = \begin{cases} 1 & \text{if } A_v \text{ depends on } T_\alpha, \\ 0 & \text{otherwise}, \end{cases}$$

for any $v = 1, \ldots, m$ and $\alpha = 0, \ldots, n$. For example, considering the example introduced in Section 3, with $m = 2, n = 1, A_1(t_1) = R(T_1)$ and $A_2(t_2) = \Phi(T_1)$, we have the following dependence matrix,

$$M_v = \begin{pmatrix} 0 & 1 \\ 0 & 1 \end{pmatrix} \iff A_1(t_1) \iff A_2(t_2).$$

The efficiency of the algorithms for the basic algebra of a series depends on the way is coded. An overcoded structure that makes good use of memory generally requires complex algorithms, which increase the computational cost in terms of time. On the other hand, an undercoded computational representation of the series generates simple algorithms, because the location of all the coefficients can be obtained directly [17]. However, this scheme presents the inconvenience of being very wasteful in the memory resources required for the storage of the series. As pointed out in [16], most of the operations involving a series are based on navigating and searching through the structure that represents the series.

Let us introduce now the red–black tree structure. A red–black tree is a special type of tree, where each node has a color attribute, the value of which is either red or black. In addition to the ordinary requirements imposed on binary search trees, the following additional requirements of any valid red-black tree apply: A node is either red or black. The root is black. All leaves are black, even when the parent is black. Both children of every red node are black. Every simple path from a node to a descendant leaf contains the same number of black nodes. A critical property of red-black trees is enforced by these constraints: the longest path from the root to a leaf is no more than twice as long as the shortest path from the root to a leaf in that tree. The result is that the tree is roughly balanced. Because operations such as inserting, deleting, and finding values requires worst-case time proportional to the height of the tree, this fact makes the red–black tree be a efficient. For instance, the search–time results to be $O(\log n)$. 

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In this figure, it is assumed that one index fails. The data associated to each node of the tree is a real number representing the coefficient of the corresponding term (\(C_i\)), and the key of each node is given by the set \(i_1, \sigma_{i_1,2}, \sigma_{i_1,3}, \ldots, i_m, \sigma_{i_m,2}, \sigma_{i_m,1}\). In Figure 2, we show the tree structure in which a polynomial is stored. The polynomial we show in Figure 2 has seven terms, with keys \(v_1 < v_2 < \ldots < v_7\). Each key corresponds to a vector of integer numbers of the form

\[
v = (v_1, \sigma_{v_1,2}, \sigma_{v_1,3}, \ldots, v_m, \sigma_{v_m,2}, \sigma_{v_m,1})
\]

If we store the key of a term in a vector structure, the complexity of the comparison of the keys is \(O(m)\). We can reduce this complexity by storing keys in red–black trees. For each term of a Poisson series, we store vectors \((v, i_1, \sigma_{i_1,2}, \sigma_{i_1,3})\). Here, \(v\) represents the number of the corresponding variable. Thus, the complexity of comparison between terms is reduced from \(O(m)\) to \(O(\log_2(m))\) in the worst case scenario. If the keys associated to different terms have different size, that means that both terms are not equal and cannot be collected. This fact helps also to reduce the computation time. Moreover, it is not necessary to compare the entire key in case one index fails.

**Figure 1.** A red–black tree represents a Poisson series. Here, \(v_1 = (v_1,0, v_{1,1}, \ldots, v_{1,m}), v_2 = (v_2,0, v_{2,1}, \ldots, v_{2,m})\) and \(v_3 = (v_3,0, v_{3,1}, \ldots, v_{3,m})\), where \(v_1 < v_2 < v_3\). Thus, the most adequate structure for storing a modified Poisson series is a red–black tree for the storage of the trigonometric part, where each node of the structure stores a key, given by the vector of integers \(v = (v_0, v_1, \ldots, v_m)\) representing the angle of the trigonometric part. Each of these nodes is linked to two red–black trees for the storage of the polynomials \(P_v\) and \(Q_v\), respectively.

In Figure 1, we show the structure used to represent the modified Poisson series given by

\[
S = \sum_{i=1}^{m} P_{vi}(A_1(t_1), \ldots, A_m(t_m)) \sin(v_{1,0}T_0 + \sum_{i=1}^{m} v_{1,i}A_i(t_i)) + P_{v2}(A_1(t_1), \ldots, A_m(t_m)) \cos(v_{2,0}T_0 + \sum_{i=1}^{m} v_{2,i}A_i(t_i)) + P_{v3}(A_1(t_1), \ldots, A_m(t_m)) \sin(v_{3,0}T_0 + \sum_{i=1}^{m} v_{3,i}A_i(t_i)) + P_{v4}(A_1(t_1), \ldots, A_m(t_m)) \cos(v_{4,0}T_0 + \sum_{i=1}^{m} v_{4,i}A_i(t_i)) + P_{v5}(A_1(t_1), \ldots, A_m(t_m)) \sin(v_{5,0}T_0 + \sum_{i=1}^{m} v_{5,i}A_i(t_i)) + P_{v6}(A_1(t_1), \ldots, A_m(t_m)) \cos(v_{6,0}T_0 + \sum_{i=1}^{m} v_{6,i}A_i(t_i)) + P_{v7}(A_1(t_1), \ldots, A_m(t_m)) \sin(v_{7,0}T_0 + \sum_{i=1}^{m} v_{7,i}A_i(t_i)).
\]

In this figure, it is assumed that \(v_1 < v_2 < v_3\).

Now, we will concentrate on the way the polynomial parts \(P_v\) and \(Q_v\) are stored. As explained earlier, each node of the Poisson series structure is linked with two red–black trees that represent \(P_v\) and \(Q_v\), respectively. Let us now consider one of these two polynomials,

\[
P_v(A_1(t_1), \ldots, A_m(t_m)) = \sum_{i \in \mathbb{C}} C \sigma_i A_i^0(t_1) D_{\sigma_1} A_1(t_1) \times A_i^1(t_2) D_{\sigma_2} A_2(t_2) \times \cdots \times A_i^m(t_m) D_{\sigma_m} A_m(t_m).
\]
Figure 2. Red–black tree for representing a polynomial with keys $v_1 < v_2 < v_3 < v_4 < v_5 < v_6 < v_7$. Each node of the tree contains also the value of the coefficient of the term ($C_i$).

Figure 3. Red–black tree structure for the storage of the modified Poisson series given in equation (24).

Figure 4. Representation of the key $K = A_1(t_1)A_2(t_2)A_3(t_3)D_{1,0}A_3(t_3)$ in a red–black tree.

Thus, from a computational point of view, a modified Poisson series will be represented by a red–black tree where each node is linked to two red–black trees with keys stored in red–black trees. In Figure 3, we show the representation of the Poisson series

$$S = (A_1(t_1) - A_1^2(t_1)D_{1,0}A_2(t_2) + A_1^2(t_1)A_2(t_2)D_{1,1}A_2(t_2)) \sin(T_0 + A_1(t_1)) -$$

$$+ (A_2(t_2)) \cos(T_0 + A_1(t_1)) +$$

$$+ (A_1^2(t_1) - A_1^2(t_1)D_{1,0}A_2(t_2) + A_1^2(t_1)A_2^2(t_2)D_{1,1}A_2(t_2)) \sin(T_0 + 2A_1(t_1)) +$$

$$+ (A_2(t_2) - A_1(t_1)A_2(t_2) + A_1(t_1)) \cos(T_0 + 2A_1(t_1)) +$$

$$+ (A_1(t_1)D_{1,1}A_2(t_2)) \sin(3T_0 + A_1(t_1) - A_2(t_2)) +$$

$$+ (A_2(t_2) - A_1(t_1)A_2^2(t_2) + A_1^2(t_1)) \cos(3T_0 + A_1(t_1) - A_2(t_2)), \tag{24}$$

just to clarify the way red–black trees are used to store a Poisson series.

We also illustrate the way the key is coded in a red–black tree structure. In Figure 4, we show the representation of the key

$$K = A_1^2(t_1)A_2^2(t_2)A_3^2(t_3)D_{1,0}A_3(t_3),$$
considering that \( m = 3 \). It is convenient to say that we need to store the complete key in the tree structure, given by

\[
K = A_1^2(t_1)D_{-1,-1}A_1(t_1)A_2^2(t_2)D_{-1,-1}A_2(t_2)A_3(t_3)D_{1,0}A_3(t_3).
\]

Here, we have the following vectors to be stored in the red–black tree structure: \((1, 2, -1, -1)\), \((2, 4, -1, -1)\), and \((3, 1, 1, 0)\).

6. Conclusions

In this paper, we analyze the method of multiple scales from the point of view of the symbolic computation to find the mathematical object that is needed to implement a general algorithm to apply the method of multiple scales in its standard version. As a consequence of the analysis of the method, we see that a modification of a Poisson series is enough to handle the expressions involved in the application of the method of multiple scales. In Section 4, we study the way the structure of a Poisson series must be adapted in order to implement the method as a symbolic algorithm. Finally, we deal with the most adequate structure for storing these series, a red–black tree for the storage of the trigonometric part, where each node of the structure stores a key, given by the vector of integers representing the angle of the trigonometric part. Each of these nodes is linked to two red–black trees for the storage of the polynomials \( P_n \) and \( Q_m \), respectively. Moreover, we propose storing the key of a term in a vector structure, so the complexity of the comparison is reduced from \( O(m) \) to \( O(\log_2 m) \) in the worst case scenario.

References

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