# SERIES EXPANSIONS FOR THE EVOLUTION OF MECHANICAL CONTROL SYSTEMS* 

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#### Abstract

This paper presents a series expansion that describes the evolution of a mechanical system starting at rest and subject to a time-varying external force. Mechanical systems are presented as second-order systems on a configuration manifold via the notion of affine connections. The series expansion is derived by exploiting the homogeneity property of mechanical systems and the variations of constant formula. A convergence analysis is obtained using some analytic functions and combinatorial analysys results. This expansion provides a rigorous mean of analyzing locomotion gaits in robotics, and lays the foundation for the design of motion control algorithms for a large class of underactuated mechanical systems.


Key words. series expansions, control of mechanical systems, nonlinear controllability

AMS subject classifications. 30B99, 70Q05, 93B05, 93B29

1. Introduction. The general purpose of this work is to develop innovative and powerful control and analysis method for underactuated mechanical control systems. This paper introduces a series expansion that characterize the evolution of a mechanical system starting at rest and subject to an open loop time-varying force. This tool should prove useful in the study of robotic locomotion and in the design of motion control algorithms.
1.1. Series expansions and their control applications. Original works on perturbation methods and series expansions in mechanics go back to Poincarè and Lagrange. Magnus [34] describes the evolution of systems on a Lie group. Chen [15], Fliess [18], and Sussmann [43] develop a general framework to describe the evolution of a nonlinear system via the so-called Chen-Fliess series and its factorization. Related work on the "chronological calculus" formalism was developed by Agračhev and Gamkrelidze [2].

Within the context of modern geometric control theory, series expansions play a key role in the study of nonlinear controllability. Small-time local controllability was studied for example by Sussmann [42, 44], Agračhev and Gamkrelidze [3], and Kawski [23, 25]. Controllability along trajectories is investigated by Bianchini and Stefani in [8]. Finally, the work by Lewis and Murray [33] on configuration controllability for mechanical control systems is very related to this work.

Motion planning problems provide a second important use of series expansions. A rich literature is available on the motion planning problem for kinematic systems, that is systems without drift. Numerous approaches include algorithms for chained systems by Murray and Sastry [37], for systems on Lie groups by Leonard and Krishnaprasad [31] and Kolmanovsky and McClamroch [27], and the very general solution proposed by Lafferriere and Sussmann [29]. These works rely on the following observation: an explicit expression for the "input history to final displacement" map

[^0]simplifies dramatically the two-point boundary value problem that defines the motion planning task. In other words, whenever an explicit expression (provided by a series expansion) for the evolution of the control system is available, the two-point boundary value problem is reduced to a low dimensional nonlinear program. Accordingly, motion control algorithms are designed by inverting this "input history to final displacement" map.

Finally, series expansions and the techniques developed in this paper have potential relevance in several areas including averaging and vibrational stabilization [6, 10], high order variations for use in optimal control [26], digital multirate sampling of nonlinear systems [19] and model reduction [20].

Series expansions that specifically exploit the structure of mechanical systems have so far not been computed. However, some preliminary progress in this direction has been obtained by Bullo, Leonard and Lewis [12, 14] via a perturbation analysis. Under the assumption of small amplitude forcing, the authors compute the initial terms of a Taylor series describing the forced evolution. The results are then found to be in agreement with the controllability analysis in [33]. A different but related research direction has focused on open loop vibrational control and the recent progress we described in [10] is related to this paper.
1.2. Summary of results. The main contribution of the paper is a series that describes the evolution of a forced mechanical system starting from rest. Mechanical systems are characterized as second-order systems on a configuration manifold using the theory of affine connections. By exploiting the problem's structure, the system's evolution is described as a flow on the configuration space ( $n$-dimensional) instead of a series on the full phase space ( $2 n$-dimensional).

The treatment relies on some differential geometric tools to describe the homogeneity properties of nonlinear mechanical systems and the variations of constants formula; see [2] and [24]. The homogeneous structure of nonlinear mechanical systems leads to a recursive procedure to compute the forced solution to a mechanical system. The terms in the series are computed recursively via time integrals and certain Lie brackets called symmetric products [33].

The series is guaranteed to convergence in a strong sense for small amplitude inputs and bounded final time. The convergence analysis is sophisticated and relies on various concepts from complex and combinatorial analysis. Following the analysis by Agračhev and Gamkrelidze in [2, Proposition 2.1], a bound is computed for every term of the series so that a notion of order is established. However, as opposed to [2], only a recursive expression for the series terms is available and this much complicates the treatment. The key idea is to obtain a recursive bound not only on the terms of the expansion but also on their partial derivative.

The series expansion can be computed in simplified fashion in two settings. For simple Hamiltonian systems with integrable forces, the main theorem can be interpreted as a statement on gradient and Hamilton flows: the flow of a Hamiltonian system forced from rest can be written as a (time-varying) gradient flow. For invariant systems on groups, the series can be computed via algebraic manipulations (no differentiations). In other words, the computations are performed on the corresponding Lie algebra and the theorem reduces to a statement on the flow of polynomial control systems. These results agree and supersede the preliminary results in in [12, 14].

Finally, some numerical simulations of a three degree of freedom robotic manipulator are performed. Truncating the series expansion at increasingly higher order, various approximations are obtained and their accuracy is illustrated via some nu-
merical data.
1.3. Organization. The paper is organized as follows. In Section 2 we present the model and the homogeneity properties of a large class of mechanical control systems. Most ideas are common in the literature, some are not. In Section 3 we present the main result of the paper, that is, a convergent series describing the evolution of a forced mechanical system. Section 4 contains some applications and extensions, including the simple Hamiltonian and the invariant system settings, as well as some simulations. We present our conclusions in Section 5.
2. Some Geometric and Analytical Properties of Mechanical Systems. We present a geometric definition of mechanical control systems, study their homogeneous properties, and provide bounds using analytic function theory.
2.1. Natural objects on manifolds. We review some basic definitions to fix some notation, see [1]. All the objects we consider are smooth in the sense of analytic. Let $Q$ be a finite dimensional, Hausdorff, second countable manifold, $q$ be a point on it, $v_{q}$ be a point on $T Q, I \subset \mathbb{R}$ be a real interval and $\gamma: I \rightarrow Q$ be a curve on $Q$. We let $0_{q}$ denote the zero velocity tangent vector at on the tangent space $T_{q} Q$. Let $\pi: T Q \rightarrow Q$ denote the usual projection on the tangent bundle, that is, $\pi\left(v_{q}\right)=q$. On the manifold $Q$, we will define scalar functions $q \mapsto f(q) \in \mathbb{R}$ and vector fields $q \mapsto X(q) \in T_{q} Q$. Lie derivatives of functions and Lie brackets of vector fields are denoted by

$$
\mathscr{L}_{X} f, \quad \text { and } \quad \mathscr{L}_{X} Y=[X, Y] .
$$

2.2. Variation of constants formula in geometric terms. This section presents a quick review of the variation of constants formula within the chronological calculus formalism introduced in [2], see also [39]. Given a vector field $Y$ and a diffeomorphism $\phi$, the pull-back of $Y$ along $\phi$, denoted $\phi^{*} Y$, is a vector field defined by

$$
\left(\phi^{*} Y\right)(q) \triangleq T_{q} \phi^{-1} \circ Y \circ \phi,
$$

where $T_{q} \phi^{-1}$ is the tangent map to $\phi^{-1}$, see [1]. In a system of local coordinate $\left(q^{1}, \ldots, q^{n}\right)$, a vector field is written as $Y(q)=Y^{i}(q) \partial / \partial q^{i}$, and the pull-back of $Y$ along $\phi$ is

$$
\left(\phi^{*} Y\right)^{i}(q)=\frac{\partial\left(\phi^{-1}\right)^{i}}{\partial q^{j}} Y^{j}(\phi(q),
$$

where the summation convention is enforced here and in what follows.
A time-varying vector field $(q, t) \mapsto X(q, t)$ gives rise to the initial value problem

$$
\dot{q}(t)=X(q, t), \quad q(0)=q_{0},
$$

and its solution at time $T$, which we refer to as the flow of $X$, is denoted by $q(T)=$ $\Phi_{0, T}^{X}\left(q_{0}\right)$. We shall usually assume time-varying quantities to be integrable with respect to time. Given a time-varying vector field $X(q, t)$, we denote its definite time integral from time 0 to time $T$ by:

$$
\begin{equation*}
\bar{X}(q, T)=\int_{0}^{T} X(q, \tau) d \tau . \tag{2.1}
\end{equation*}
$$



Fig. 2.1. The flow along $X+Y$ is written as the flow along $Y$ with initial condition $q_{q}$. The "variation" $q_{1}$ is computed via the variation of constants formula as the flow along $\left(\Phi_{0, t}^{Y}\right)^{*} X$ for time $[0, T]$ with initial condition $q_{0}$.

The integral takes place over the linear space $T_{q} Q$ at fixed $q \in Q$. This operation can be defined in two ways. Given a coordinates chart about $q$, the integral is well-defined in the coordinate system (this definition suffices for the purpose of this paper, since the analysis is local). A global coordinate-free definition is obtained providing sufficient conditions in order for $T_{q} Q$ to be a Banach space and introducing the Cauchy-Bochner integral, see [1, see the discussion at page 61].

Next, consider the initial value problem

$$
\begin{equation*}
\dot{q}(t)=X(q, t)+Y(q, t), \quad q(0)=q_{0} \tag{2.2}
\end{equation*}
$$

where $X$ and $Y$ are analytic time-varying vector fields. If we regard $X$ as a perturbation to the vector field $Y$, we can describe the flow of $X+Y$ in terms of a nominal and perturbed flow. The following relationship is referred to as the variation of constants formula and describes the perturbed flow:

$$
\begin{equation*}
\Phi_{0, t}^{X+Y}=\Phi_{0, t}^{Y} \circ \Phi_{0, t}^{\left(\Phi_{0, t}^{Y}\right)^{*} X} \tag{2.3}
\end{equation*}
$$

The result is illustrated in Figure 2.1 and proven in [2, equation (3.15)], see also [10, Appendix A.1]. The result can be alternatively stated as follows. For all $T \geq 0$, the final value $q(T)$ of the curve $q:[0, T] \rightarrow M$ solution to the initial value problem (2.2) is also the final value of the curve solution to

$$
\begin{equation*}
\dot{q}(s)=Y(q, s), \quad q(0)=z(T) \tag{2.4}
\end{equation*}
$$

where $z:[0, T] \rightarrow M$ is the solution to the initial value problem

$$
\begin{equation*}
\dot{z}(s)=\left(\left(\Phi_{0, s}^{Y}\right)^{*} X\right)(z), \quad z(0)=q_{0} \tag{2.5}
\end{equation*}
$$

The differential equation (2.5) is referred to as the "pulled back" or the "adjoint" system in [21]. If both $X$ and $Y$ are time invariant, then the classic infinitesimal Campbell-Backer-Hausdorff formula, see [22], provides a mean of computing the pullback:

$$
\left(\Phi_{0, t}^{Y}\right)^{*} X=\sum_{k=0}^{\infty} \operatorname{ad}_{Y}^{k} X \frac{t^{k}}{k!}
$$

If instead $X$ and $Y$ are time-varying, a generalized expression is, see [2]:

$$
\begin{align*}
\left(\Phi_{0, t}^{Y}\right)^{*} X(q, t)= & X(q, t) \\
& +\sum_{k=1}^{\infty} \int_{0}^{t} \ldots \int_{0}^{s_{k-1}}\left(\operatorname{ad}_{Y\left(q, s_{k}\right)} \ldots \operatorname{ad}_{Y\left(q, s_{1}\right)} X(q, t)\right) d s_{k} \ldots d s_{1} . \tag{2.6}
\end{align*}
$$

Just like in the classic Campbell-Backer-Hausdorff formula, see [45], the convergence of the series expansion in the previous equation is a delicate manner. Sufficient conditions for local convergence are given in [2, Proposition 2.1 and 3.1]. For our analysis, the following straightforward statement suffices. If all the Lie brackets $\operatorname{ad}_{Y\left(s_{k}\right)} \ldots \operatorname{ad}_{Y\left(s_{1}\right)} X$ vanish for all $k$ greater than a given $N$, then the series in equation (2.6) becomes a finite sum and it readily converges.
2.3. Affine connections. We refer to $[17,30]$ for a comprehensive treatment on affine connections and Riemannian geometry. An affine connection on $Q$ is a smooth map that assigns to a pair of vector fields $X, Y$ a vector field $\nabla_{X} Y$ such that for any function $f$ and for any third vector field $Z$ :
(i) $\nabla_{f X+Y} Z=f \nabla_{X} Z+\nabla_{Y} Z$,
(ii) $\nabla_{X}(f Y+Z)=\left(\mathscr{L}_{X} f\right) Y+f \nabla_{X} Y+\nabla_{X} Z$.

We also say that $\nabla_{X} Y$ is the covariant derivative of $Y$ with respect to $X$. Vector fields can also be covariantly differentiated along curves, and this concept will be instrumental in writing the Euler-Lagrange equations. Consider a smooth curve $\gamma$ : $[0,1] \rightarrow Q$ and a vector field along $\gamma$, that is, a map $v:[0,1] \rightarrow T Q$ such that $\pi(v(t))=$ $\gamma(t)$ for all $t \in[0,1]$. Let $V$ be a smooth vector field satisfying $V(\gamma(t))=v(t)$. The covariant derivative of the vector field $v$ along $\gamma$ is defined by

$$
\frac{D v(t)}{d t} \triangleq \nabla_{\dot{\gamma}(t)} v(t)=\left.\nabla_{\dot{\gamma}(t)} V(q)\right|_{q=\gamma(t)}
$$

It can be shown that this definition is independent of the choice of $V$. In a system of local coordinate $\left(q^{1}, \ldots, q^{n}\right)$, an affine connection is uniquely determined by its Christoffel symbols ${ }^{1} \Gamma_{i j}^{i}$

$$
\nabla_{\frac{\partial}{\partial q^{i}}}\left(\frac{\partial}{\partial q^{j}}\right)=\Gamma_{i j}^{k} \frac{\partial}{\partial q^{k}}
$$

and accordingly, the covariant derivative of a vector field is written as

$$
\nabla_{X} Y=\left(\frac{\partial Y^{i}}{\partial q^{j}} X^{j}+\Gamma_{j k}^{i} X^{j} Y^{k}\right) \frac{\partial}{\partial q^{i}} .
$$

2.4. Control systems described by affine connections. We introduce a class of control systems that is a generalization of Lagrangian control systems. This approach to modeling of vehicles and robotic manipulators is common to a number of recent works; see $[9,33,32,10]$. A control system described by an affine connection is defined by the following objects:
(i) an $n$-dimensional configuration manifold $Q$, with $q \in Q$ being the configuration of the system and $v_{q} \in T_{q} Q$ being the system's velocity,

[^1](ii) an affine connection $\nabla$ on $Q$, whose Christoffel symbols are $\left\{\Gamma_{j k}^{i}: i, j, k \in\right.$ $\{1, \ldots, n\}\}$,
(iii) a time-varying vector field $Y$ on $Q$ defining the input force.

The corresponding equations of motion are written as

$$
\begin{equation*}
\frac{D v_{q}}{d t}=Y(q, t) \tag{2.7}
\end{equation*}
$$

or equivalently in coordinates as

$$
\begin{equation*}
\dot{q}^{i}=v^{i}, \quad \dot{v}^{i}+\Gamma_{j k}^{i}(q) v^{j} v^{k}=Y^{i}(q, t) \tag{2.8}
\end{equation*}
$$

where the indices $i, j, k$ run from 1 to $n$, and where $v_{q}=v^{i} \frac{\partial}{\partial q^{i}}$. These equations are a generalized form of the Euler-Lagrange equations.

REMARK 2.1. This definition of control systems described by an affine connection provides a convenient mean of treating various classes of Lagrangian mechanical systems. For example, systems with nonholonomic constraints are described within this framework in [32]. We will treat in more details "simple Hamiltonian systems" in Section 4.2 and "invariant systems on Lie groups" in Section 4.3. A more detailed exposition in presented in [10].

The second-order system in equation (2.7) can be written as a first-order differential equation on the tangent bundle $T Q$. Using $\left\{\frac{\partial}{\partial q^{i}}, \frac{\partial}{\partial v^{i}}\right\}$ as basis for the tangent bundle to $T Q$, we define

$$
Z\left(v_{q}\right)=v^{i} \frac{\partial}{\partial q^{i}}-\Gamma(q)_{j k}^{i} v^{j} v^{k} \frac{\partial}{\partial v^{i}}, \quad \text { and } \quad Y^{\mathrm{lift}}\left(v_{q}, t\right)=Y^{i}(q, t) \frac{\partial}{\partial v^{i}},
$$

so that the control system is rewritten as

$$
\begin{equation*}
\dot{v}_{q}=Z\left(v_{q}\right)+Y^{\mathrm{lift}}\left(v_{q}, t\right) \tag{2.9}
\end{equation*}
$$

We refer to $[33,30]$ for coordinate independent definitions of the lifting operation $Y \rightarrow Y^{\text {lift }}$ and of the drift vector field $Z$.
2.5. Homogeneity and Lie algebraic structure. One fundamental structure of the control system in equation (2.7) is the polynomial dependence of the vector fields $Z$ and $Y^{\text {lift }}$ on the velocity variables $v_{q}^{i}$. This structure is reflected in the Lie brackets computations involving $Z$ and $Y^{\text {lift }}$; see related ideas in [41, 10].

We here rely on the notion of geometric homogeneity ${ }^{2}$ as described in [24]. Given two vector fields $X$ and $X_{E}$, we say that the vector field $X$ is homogeneous with degree $m$ with respect to $X_{E}$ if

$$
\left[X_{E}, X\right]=m X
$$

For control systems described by an affine connection, we introduce the Liouville vector field on $T Q$, see [ 7 , pages 19 and 29], as

$$
L\left(v_{q}\right)=v^{i} \frac{\partial}{\partial v^{i}}
$$

[^2]

Fig. 2.2. Table of Lie brackets between the drift vector field $Z$ and the input vector field $Y^{l i f t}$. The $(i, j)$ th position contains Lie brackets with $i$ copies of $Y^{\text {lift }}$ and $j$ copies of $Z$. The corresponding homogeneous degree is $j-i$. All Lie brackets to the right of $\mathcal{P}_{-1}$ exactly vanish. All Lie brackets to the left of $\mathcal{P}_{-1}$ vanish when evaluated at $v_{q}=0_{q}$.
where we recall $v_{q}=v^{i} \frac{\partial}{\partial q^{i}}$. The key mathematical relationships between vector fields on $T Q$ are

$$
[L, Z]=(+1) Z, \quad \text { and } \quad\left[L, Y^{\mathrm{lift}}\right]=(-1) Y^{\mathrm{lift}}
$$

Hence, the vector field $Z$ is homogeneous of degree +1 , and the vector field $Y^{\text {lift }}$ is homogeneous of degree -1 with respect to the Liouville vector field. Let $\mathcal{P}_{j}$ be the set of vector fields on $T Q$ of homogeneous degree $j$, so that

$$
Z \in \mathcal{P}_{1}, \quad \text { and } \quad Y^{\text {lift }} \in \mathcal{P}_{-1}
$$

The sets $\mathcal{P}_{j}$ enjoy various interesting properties. Table 2.2 illustrates them, their proof is via direct computation, and they are listed next:
(i) $\left[\mathcal{P}_{i}, \mathcal{P}_{j}\right] \subset \mathcal{P}_{i+j}$, that is, the Lie bracket between a vector field in $\mathcal{P}_{i}$ and a vector field in $\mathcal{P}_{j}$ belongs to $\mathcal{P}_{i+j}$.
(ii) $\mathcal{P}_{k}=\{0\}$ for all $k \leq-2$,
(iii) for all $X \in \mathcal{P}_{k}$ with $k \geq 1, X\left(0_{q}\right)=0_{q}$,
(iv) every $X \in \mathcal{P}_{-1}$ is the lift of a vector field on $Q$.

It is helpful to provide an interpretation of $\mathcal{P}_{i}$ in coordinates. In a system of local coordinates, let $\mathcal{H}_{i}\left(q, v_{q}\right)$ be the set of scalar functions on $T Q=\mathbb{R}^{2 n}$, which are arbitrary functions of $q$ and which are homogeneous polynomials in $\left\{v^{1}, \ldots, v^{n}\right\}$ of degree $i . \mathcal{P}_{i}$ is the set of vector fields on $\mathbb{R}^{2 n}$ with the first $n$ components in $\mathcal{H}_{i}$ and the second $n$ components in $\mathcal{H}_{i+1}$.

Finally, it is of interest to focus on the Lie bracket $\left[Y_{b}^{\text {lift }},\left[Z, Y_{a}^{\text {lift }}\right]\right]$, where $Y_{a}, Y_{b}$ are two vector fields on $Q$. This operation will play an important role in later computations. Since this Lie bracket belongs to $\mathcal{P}_{-1}$, there must exist a vector field on $Q$, which we denote $\left\langle Y_{a}: Y_{b}\right\rangle$, such that

$$
\left\langle Y_{a}: Y_{b}\right\rangle^{\text {lift }}=\left[Y_{b}^{\text {lift }},\left[Z, Y_{a}^{\text {lift }}\right]\right] .
$$

Such a vector field is called symmetric product between $Y_{b}$ and $Y_{a}$ and a direct computation shows that it satisfies

$$
\left\langle Y_{b}: Y_{a}\right\rangle=\nabla_{Y_{a}} Y_{b}+\nabla_{Y_{b}} Y_{a}
$$

or equivalently in coordinates

$$
\left\langle Y_{b}: Y_{a}\right\rangle^{i}=\frac{\partial Y_{a}^{i}}{\partial q^{j}} Y_{b}^{j}+\frac{\partial Y_{b}^{i}}{\partial q^{j}} Y_{a}^{j}+\Gamma_{j k}^{i}\left(Y_{a}^{j} Y_{b}^{k}+Y_{a}^{k} Y_{b}^{j}\right)
$$

The adjective "symmetric" comes from the equality $\left\langle Y_{a}: Y_{b}\right\rangle=\left\langle Y_{b}: Y_{a}\right\rangle$.
2.6. Integrable flows. Here we compute solutions to a few differential equations defined by certain homogeneous vector fields. In particular, significant simplifications take place in the following two cases. First, let $(q, t) \mapsto X(q, t)$ be a time-varying vector field on $Q$, and consider the differential equation on $T Q$

$$
\begin{equation*}
\dot{v}_{q}=X^{\mathrm{lift}}\left(v_{q}, t\right) \tag{2.10}
\end{equation*}
$$

with initial condition $v_{q}(0)=v_{0} \in T_{q_{0}} Q$. It can be seen that

$$
\begin{equation*}
\Phi_{0, t}^{X^{\mathrm{lift}}}\left(v_{0}\right)=v_{0}+\int_{0}^{t} X\left(q_{0}, s\right) d s \tag{2.11}
\end{equation*}
$$

that is, in coordinates

$$
\Phi_{0, t}^{X^{\mathrm{lift}}}\left(\left[\begin{array}{l}
q_{0} \\
v_{0}
\end{array}\right]\right)=\left[\begin{array}{c}
q_{0} \\
v_{0}+\int_{0}^{t} X\left(q_{0}, s\right) d s
\end{array}\right]
$$

Next, let $X_{0} \in \mathcal{P}_{0}$ and $X_{1} \in \mathcal{P}_{1}$ and consider the differential equation

$$
\begin{equation*}
\dot{v}_{q}=X_{0}\left(v_{q}, t\right)+X_{1}\left(v_{q}, t\right), \tag{2.12}
\end{equation*}
$$

with initial condition $v_{q}(0)=0_{q_{0}} \in T_{q_{0}} Q$. Define the vector field $X_{0,1}$ on $Q$ and its flow $\zeta:[0, T] \mapsto Q$ via

$$
\begin{aligned}
X_{0,1} & =T \pi \circ X_{0} \\
\zeta(t) & =\Phi_{0, t}^{X_{0,1}}\left(q_{0}\right),
\end{aligned}
$$

where $T \pi: T T Q \rightarrow T Q$ is the tangent map to the projection map $\pi: T Q \rightarrow Q$. In coordinates, this vector field consists of the first $n$ components of the vector field $X_{0}=\left[X_{0,1}(q, t)^{\prime}, \quad X_{0,2}(q, v, t)^{\prime}\right]^{\prime}$ on $T Q$. It can be seen that

$$
\Phi_{0, t}^{X_{0}+X_{1}}\left(0_{q_{0}}\right)=0_{\zeta(t)}
$$

that is, in coordinates

$$
\Phi_{0, t}^{X_{0}+X_{1}}\left(\left[\begin{array}{c}
q_{0} \\
0
\end{array}\right]\right)=\left[\begin{array}{c}
\zeta(t) \\
0
\end{array}\right]
$$

The key observation in proving this statement is that the components of $X_{0,2}$ and $X_{1}$ are polynomials in $\left\{v^{1}, \ldots, v^{n}\right\}$ of degree at least 1 . Since the initial velocity is assumed zero, $v_{q}$ remains zero for all time.
2.7. Analyticity and bounds over complex neighborhoods. In this section we introduce a norm on the set of analytic vector fields over a compact subset of $Q$. We also provide bounds to partial derivatives of analytic functions. The bounds are not coordinate-free, i.e., they depend on the specific selection of coordinate system. Accordingly, the treatment here assumes $Q=\mathbb{R}^{n}$.

Let $q_{0}$ be a point on $\mathbb{R}^{n}$, and let $\sigma$ be a positive scalar, and define the complex $\sigma$-neighborhood of $q_{0}$ in $\mathbb{C}^{n}$ as

$$
B_{\sigma}\left(q_{0}\right)=\left\{z \in \mathbb{C}^{n}:\left\|z-q_{0}\right\|<\sigma\right\}
$$

Let $f$ be a real analytic function on $\mathbb{R}^{n}$ that admits a bounded analytic continuation over $B_{\sigma}\left(q_{0}\right)$. The norm of $f$ is defined as

$$
\|f\|_{\sigma} \triangleq \max _{z \in B_{\sigma}\left(q_{0}\right)}|f(z)|
$$

where $f$ denotes both the function over $\mathbb{R}^{n}$ and its analytic continuation. Given a time-varying vector field $(q, t) \mapsto Y(q, t)=Y_{t}(q)$, let $Y_{t}^{i}$ be its $i$ th component with respect to the usual basis on $\mathbb{R}^{n}$. Assuming $t \in[0, T]$, and assuming that every component function $Y_{t}^{i}$ is analytic over $B_{\sigma}\left(q_{0}\right)$, we define the norm of $Y$ as

$$
\|Y\|_{\sigma, T} \triangleq \max _{t \in[0, T]} \max _{i \in\{1, \ldots, n\}}\left\|Y_{t}^{i}\right\|_{\sigma}
$$

In what follows, we will often simplify notation by neglecting the subscript $T$ in the norm of a time-varying vector field. Finally, given an affine connection $\nabla$ with Christoffel symbols $\left\{\Gamma_{j k}^{i}: i, j, k \in\{1, \ldots, n\}\right\}$, introduce the notation:

$$
\|\Gamma\|_{\sigma} \triangleq \max _{i j k}\left\|\Gamma_{j k}^{i}\right\|_{\sigma}
$$

Next, we examine the norm of partial derivatives of these objects. Recall that the Cauchy integral representation of analytic functions leads to bounds on high order derivatives of analytic functions in terms of the norm of the functions themselves, see the so-called Cauchy estimates in [28, Section 2.3] and [38]. Let $\left(i_{1}, \ldots, i_{m}\right)$ be a collection of integers belonging to $\{1, \ldots, n\}$, and let $\sigma^{\prime}$ be a positive real strictly less than $\sigma$. It is known that

$$
\left\|\partial^{m} f\right\|_{\sigma^{\prime}} \triangleq \max _{i_{1}, \ldots, i_{m}}\left\|\frac{\partial^{m} f}{\partial q_{i_{1}} \cdots \partial q_{i_{m}}}\right\|_{\sigma^{\prime}} \leq m!\delta^{m}\|f\|_{\sigma}
$$

where $\delta=n /\left(\sigma-\sigma^{\prime}\right)$. The quantity $\partial^{m} f / \partial q_{i_{1}} \cdots \partial q_{i_{m}}$ is a real function; it is bounded by bounding its analytic continuation over $B_{\sigma}\left(q_{0}\right)$. Similarly, for vector fields

$$
\left\|\partial^{m} Y\right\|_{\sigma^{\prime}} \triangleq \max _{t \in[0, T] i, i_{1}, \ldots, i_{m}} \max \left\|\frac{\partial^{m} Y_{t}^{i}}{\partial q_{i_{1}} \cdots \partial q_{i_{m}}}\right\|_{\sigma^{\prime}} \leq m!\delta^{m}\|Y\|_{\sigma}
$$

and for the Christoffel symbols

$$
\left\|\partial^{m} \Gamma\right\|_{\sigma^{\prime}} \triangleq \max _{i, j, k, i_{1}, \ldots, i_{m}}\left\|\frac{\partial^{m} \Gamma_{j k}^{i}}{\partial q_{i_{1}} \cdots \partial q_{i_{m}}}\right\|_{\sigma^{\prime}} \leq m!\delta^{m}\|\Gamma\|_{\sigma}
$$

3. A Series Expansion for Mechanical Control Systems. This section describes first a preliminary bound, then the main result of the paper, that is, a series expansion describing the evolution of a forced control system starting at rest.

Problem 3.1. Assume the functions $q \mapsto \Gamma_{j k}^{i}(q)$ and the vector field $(q, t) \mapsto$ $Y(q, t)$ analytic in $q \in Q$, and integrable in $t \in[0, T]$, for some positive time $T$. Let $\gamma:[0, T] \mapsto Q$ be the solution to the differential equation (2.7) with initial condition $\dot{\gamma}(0)=0_{q_{0}}$. Characterize $\gamma$ as a series expansion containing iterated symmetric products and time integrals of $Y$.

We start with a conservative bound.
Lemma 3.2 (Bound on evolution). Consider the system as described in Problem 3.1. Select a coordinate system about the point $q_{0} \in Q$ and let $\sigma$ be a positive constant. A sufficient condition for $\gamma([0, T])$ to be a subset of $B_{\sigma}\left(q_{0}\right)$ is that

$$
\begin{equation*}
\|Y\|_{\sigma} T^{2}<\frac{\eta^{2}\left(\sigma n^{2}\|\Gamma\|_{\sigma}\right)}{n^{2}\|\Gamma\|_{\sigma}} \tag{3.1}
\end{equation*}
$$

where the function $\eta: x \in \mathbb{R}_{+} \rightarrow[0, \pi / 2]$ is the unique solution to $\eta \tan (\eta)=x$.
Proof. Let $T_{0}<T$ be the smallest time at which the solution $\gamma$ reaches the distance $\left\|\gamma\left(T_{0}\right)-q_{0}\right\|=\sigma$. If the solution never reaches this distance, then $\gamma([0, T])$ is obviously a subset of $B_{\sigma}\left(q_{0}\right)$. Since $\gamma\left(\left[0, T_{0}\right]\right) \subset B_{\sigma}\left(q_{0}\right)$, for all $t \in\left[0, T_{0}\right]$ we have the bound $\|\dot{\gamma}(t)\| \leq y(t)$, where

$$
\dot{y}=n^{2}\|\Gamma\|_{\sigma} y^{2}+\|Y\|_{\sigma}, \quad y(0)=0
$$

The solution to this initial value problem is

$$
y(t)=\sqrt{\frac{\|Y\|_{\sigma}}{n^{2}\|\Gamma\|_{\sigma}}} \tan \left(\sqrt{\|Y\|_{\sigma} n^{2}\|\Gamma\|_{\sigma}} t\right)
$$

Straightforward manipulations show that the condition in equation (3.1) is equivalent to $T y(T)<\sigma$. But since $y$ is a monotone function, also $T_{0} y\left(T_{0}\right)<\sigma$. Note that $\left\|\gamma(0)-q_{0}\right\|=0$ and

$$
\frac{d}{d t}\left\|\gamma(t)-q_{0}\right\| \leq\|\dot{\gamma}\| \leq y(t)<\sigma / T_{0}
$$

for all $t \in\left[0, T_{0}\right]$. Therefore $\left\|\gamma\left(T_{0}\right)-q_{0}\right\|<T_{0} \sigma / T_{0}$ and the contradiction is now immediate.

We are now ready to present the main theorem.
Theorem 3.3 (Evolution of a forced mechanical system starting at rest). Consider the system as described in Problem 3.1. Define recursively the time-varying vector fields $V_{k}$ :

$$
\begin{align*}
& V_{1}(q, t)=\int_{0}^{t} Y(q, s) d s  \tag{3.2}\\
& V_{k}(q, t)=-\frac{1}{2} \sum_{j=1}^{k-1} \int_{0}^{t}\left\langle V_{j}(q, s): V_{k-j}(q, s)\right\rangle d s, \quad k \geq 2 \tag{3.3}
\end{align*}
$$

Select a coordinate system about the point $q_{0} \in Q$, let $\sigma>\sigma^{\prime}$ be two positive constants, and assume that

$$
\begin{equation*}
\|Y\|_{\sigma} T^{2}<L \triangleq \min \left\{\frac{\sigma-\sigma^{\prime}}{2^{4} n^{2}(n+1)}, \frac{1}{2^{4} n(n+1)\|\Gamma\|_{\sigma}}, \frac{\eta^{2}\left(\sigma^{\prime} n^{2}\|\Gamma\|_{\sigma^{\prime}}\right)}{n^{2}\|\Gamma\|_{\sigma^{\prime}}}\right\} \tag{3.4}
\end{equation*}
$$

Then the solution $\gamma:[0, T] \rightarrow Q$ satisfies

$$
\begin{equation*}
\dot{\gamma}(t)=\sum_{k=1}^{+\infty} V_{k}(\gamma(t), t) \tag{3.5}
\end{equation*}
$$

where $V_{k}$ satisfies the bound

$$
\begin{equation*}
\left\|V_{k}\right\|_{\sigma^{\prime}} \leq L^{1-k}\|Y\|_{\sigma}^{k} t^{2 k-1} \tag{3.6}
\end{equation*}
$$

and the series $(q, t) \mapsto \sum_{k=1}^{\infty} V_{k}(q, t)$ converges absolutely and uniformly for $q \in$ $B_{\sigma^{\prime}}\left(q_{0}\right)$ and for $t \in[0, T]$.

A few comments on the various steps of the proof are appropriate. First, we investigate how to write the flow of a mechanical control system as the composition of more elementary flows. Two observations play a key role: the homogeneity of system (2.7) renders the computations tractable, the simplifying procedure can be easily repeated giving rise to an iterative procedure. Second, we prove absolute and uniform convergence of the series expansion resulting from the first formal part of the proof. The proof of the bounds is inspired by the treatment in [2, Proposition 2.1], but it is considerably more complicated here by the recursive nature of the series expansion. Once the series is formally derived and it is proven to be convergent, a limiting argument leads to the final statement in equation (3.5).

## Proof.

Part I: Here we write the solution to equation (2.7) as composition of the flow of two separate vector fields, one of which is defined recursively.

Let $k$ be a strictly positive integer, let $X_{k}, Y_{k}, W_{k}$ be time-varying vector fields on $Q$, and let $v_{q, k}$ be a smooth curve on $T Q$ that satisfies the differential equation

$$
\begin{align*}
\dot{v}_{q, k} & =\left(Z+\left[X_{k}^{\mathrm{lift}}, Z\right]+Y_{k}^{\mathrm{lift}}+W_{k}^{\mathrm{lift}}\right)\left(v_{q, k}, t\right)  \tag{3.7}\\
v_{q, k}(0) & =0_{q_{0}}
\end{align*}
$$

The mechanical system in equation (2.7) corresponds to setting $k=1, X_{1}=W_{1}=0$, $Y_{1}=Y(q, t)$, and accordingly $\dot{\gamma}(t)=v_{q, 1}(t)$. Using the formula in equations (2.4) and (2.5) discussed in Section 2.2, we set

$$
\begin{equation*}
v_{q, k}(t)=\Phi_{0, t}^{Y_{k}^{\text {lift }}}\left(v_{q, k+1}(t)\right) \tag{3.8}
\end{equation*}
$$

and

$$
\begin{align*}
\dot{v}_{q, k+1} & =\left(\left(\Phi_{0, t}^{Y_{k}^{\mathrm{lift}}}\right)^{*}\left(Z+\left[X_{k}^{\mathrm{lift}}, Z\right]+W_{k}^{\mathrm{lift}}\right)\right)\left(v_{q, k+1}\right)  \tag{3.9}\\
v_{q, k+1}(0) & =0_{q_{0}}
\end{align*}
$$

where we compute the pull-back along the flow by means of the infinite series in equation (2.6). Remarkably, this series reduces to a finite sum. From the discussion in Section 2.5 on the Lie algebraic structure of the various vector fields, we have

$$
\begin{gathered}
\operatorname{ad}_{Y_{k}^{\text {lift }}}^{m+2} Z=0 \\
\operatorname{ad}_{Y_{k}^{\text {lift }}}^{m+1}\left[X_{k}^{\text {lift }}, Z\right]=0, \quad \operatorname{ad}_{Y_{k}^{\text {lift }}}^{m} W_{k}^{\text {lift }}=0
\end{gathered}
$$

for all $m \geq 1$. With a little book-keeping we can exploit these equalities and compute

$$
\begin{aligned}
\left(\Phi_{0, t}^{Y_{k}^{\text {lift }}}\right)^{*}(Z+ & {\left.\left[X_{k}^{\text {lift }}, Z\right]+W_{k}^{\text {lift }}\right) } \\
=Z+ & {\left[X_{k}^{\text {lift }}, Z\right]+W_{k}^{\text {lift }}+\int_{0}^{t}\left[Y_{k}^{\text {lift }}(s),\left(Z+\left[X_{k}^{\text {lift }}, Z\right]\right)\right] d s } \\
& +\int_{0}^{t} \int_{0}^{s_{1}}\left[Y_{k}^{\text {lift }}\left(s_{2}\right),\left[Y_{k}^{\text {lift }}\left(s_{1}\right), Z\right]\right] d s_{2} d s_{1} \\
=Z+ & {\left[X_{k}^{\text {lift }}+\bar{Y}_{k}^{\text {lift }}, Z\right]+\left[\bar{Y}_{k}^{\text {lift }}(s),\left[X_{k}^{\text {lift }}, Z\right]\right]+W_{k}^{\text {lift }} } \\
& +\int_{0}^{t} \int_{0}^{s_{1}}\left[Y_{k}^{\text {lift }}\left(s_{2}\right),\left[Y_{k}^{\text {lift }}\left(s_{1}\right), Z\right]\right] d s_{2} d s_{1} \\
=Z & +\left[X_{k}^{\text {lift }}+\bar{Y}_{k}^{\text {lift }}, Z\right]-\left\langle\bar{Y}_{k}: X_{k}\right\rangle^{\text {lift }}+W_{k}^{\text {lift }}-\frac{1}{2}\left\langle\bar{Y}_{k}: \bar{Y}_{k}\right\rangle^{\text {lift }}
\end{aligned}
$$

where we have used the ${ }^{-}$notation introduced in equation (2.1). The last equality also relies on

$$
\int_{0}^{t} \int_{0}^{s_{1}}\left[Y_{k}^{\mathrm{lift}}\left(s_{2}\right),\left[Y_{k}^{\mathrm{lift}}\left(s_{1}\right), Z\right]\right] d s_{2} d s_{1}=-\frac{1}{2}\left\langle\bar{Y}_{k}: \bar{Y}_{k}\right\rangle^{\mathrm{lift}}
$$

which follows from an integration by part and the symmetry of the symmetric product. Remarkably, the differential equation describing the evolution of $v_{k+1}(t)$ is of the same form as equation (3.7) describing the evolution of $v_{q, k}(t)$, where

$$
\begin{aligned}
X_{k+1} & =X_{k}+\bar{Y}_{k} \\
Y_{k+1}+W_{k+1} & =-\left\langle\bar{Y}_{k} \quad: \quad X_{k}+\frac{1}{2} \bar{Y}_{k}\right\rangle+W_{k}
\end{aligned}
$$

The vector field $X_{k}$ can be computed and substituted in as:

$$
\begin{align*}
X_{k} & =\sum_{j=1}^{k-1} \bar{Y}_{j} \\
Y_{k+1}+W_{k+1} & =-\left\langle\bar{Y}_{k}: \sum_{j=1}^{k-1} \bar{Y}_{j}+\frac{1}{2} \bar{Y}_{k}\right\rangle+W_{k} \tag{3.10}
\end{align*}
$$

Notice that the quantities $Y_{k}$ and $W_{k}$ are not yet uniquely determined. Equation (3.10) is verified for all $k$, if and only if for all $m$ :

$$
\begin{equation*}
\left(Y_{2}+Y_{3}+\ldots+Y_{m+1}\right)+W_{m+1}=-\sum_{k=1}^{m}\left\langle\bar{Y}_{k}: \sum_{j=1}^{k-1} \bar{Y}_{j}+\frac{1}{2} \bar{Y}_{k}\right\rangle \tag{3.11}
\end{equation*}
$$

where we used $W_{1}=0$. Some further manipulation leads to:

$$
\begin{aligned}
\sum_{k=1}^{m}\left\langle\bar{Y}_{k}: \sum_{j=1}^{k-1} \bar{Y}_{j}+\frac{1}{2} \bar{Y}_{k}\right\rangle & =\sum_{k=1}^{m} \sum_{j=1}^{k-1}\left\langle\bar{Y}_{k}: \bar{Y}_{j}\right\rangle+\frac{1}{2} \sum_{k=1}^{m}\left\langle\bar{Y}_{k}: \bar{Y}_{k}\right\rangle \\
& =\frac{1}{2} \sum_{j, k=1, j \neq k}^{m}\left\langle\bar{Y}_{k}: \bar{Y}_{j}\right\rangle+\frac{1}{2} \sum_{k=1}^{m}\left\langle\bar{Y}_{k}: \bar{Y}_{k}\right\rangle \\
& =\frac{1}{2} \sum_{j, k=1}^{m}\left\langle\bar{Y}_{k}: \bar{Y}_{j}\right\rangle
\end{aligned}
$$

A selection of $\left\{Y_{i}: i \in\{1, \ldots, m\}\right\}$, and $W_{m+1}$ that satisfies equation (3.11) is

$$
\begin{align*}
Y_{i} & =-\frac{1}{2} \sum_{j, k=1, j+k=i}^{m}\left\langle\bar{Y}_{k}: \bar{Y}_{j}\right\rangle=-\frac{1}{2} \sum_{j=1}^{i-1}\left\langle\bar{Y}_{j}: \bar{Y}_{i-j}\right\rangle  \tag{3.12}\\
W_{m+1} & =-\frac{1}{2} \sum_{j, k=1, j+k>m}^{m}\left\langle\bar{Y}_{k}: \bar{Y}_{j}\right\rangle .
\end{align*}
$$

Note that equation (3.12) is a well defined recursive relationship, and note that the recursive definition of $V_{k}$ in equation (3.3) and (3.2) inside the theorem statement corresponds to setting $V_{k}(q, t)=\bar{Y}_{k}(q, t)$. The iteration procedure proves that, for any $k \geq 2$, the solution to the original mechanical system $\dot{\gamma}=v_{q, 1}:[0, T] \mapsto T Q$ satisfies

$$
\dot{\gamma}(t)=\left(\Phi_{0, t}^{Y_{1}^{\text {lift }}} \circ \Phi_{0, t}^{Y_{2}^{\text {lift }}} \circ \ldots \circ \Phi_{0, t}^{Y_{k-1}^{\text {lift }}}\right)\left(v_{q, k}(t)\right),
$$

where $v_{q, k}:[0, T] \mapsto T Q$ is the solution to equation (3.7). The flow $\dot{\gamma}$ is now written as the composition of $k$ flows and a first simplification is immediate. For all integers $i, j$ and for all times $s_{1}, s_{2}$ the vector fields $Y_{i}^{\text {lift }}$ and $Y_{j}^{\text {lift }}$ commute, that is

$$
\left[Y_{i}^{\mathrm{lift}}\left(v_{q}, s_{1}\right), Y_{j}^{\mathrm{lift}}\left(v_{q}, s_{2}\right)\right]=0
$$

so that $\gamma$ is the solution to

$$
\begin{equation*}
\dot{\gamma}(t)=\Phi_{0, t}^{\sum_{j=1}^{k-1} Y_{j}^{1 \mathrm{ift}}}\left(v_{q, k}(t)\right) . \tag{3.13}
\end{equation*}
$$

A second simplification is also straightforward. The vector field in equation (3.13) is homogeneous of degree 0 , i.e., it is in the form of equation (2.10) in Section 2.6. According to the result in equation (2.11) we have for all $t \in[0, T]$

$$
\begin{equation*}
\dot{\gamma}(t)=v_{q, k}(t)+\sum_{j=1}^{k-1} \bar{Y}_{j}\left(\pi\left(v_{q, k}(t)\right), t\right) \tag{3.14}
\end{equation*}
$$

where the sequence of vector fields $Y_{j}$ is defined via equation (3.12) and where the curve $v_{q, k}:[0, T] \mapsto T Q$ is the solution to

$$
\begin{align*}
\frac{d v_{q, k}}{d t} & =\left(Z+\left[\sum_{j=1}^{k-1} \bar{Y}_{j}^{\mathrm{lift}}, Z\right]+Y_{k}^{\mathrm{lift}}+W_{k}^{\mathrm{lift}}\right)\left(v_{q, k}, t\right)  \tag{3.15}\\
v_{q, k}(0) & =0_{q_{0}} .
\end{align*}
$$

Part II: Here we show absolute and uniform convergence of the series $\sum_{k=1}^{\infty} Y_{k}(q, t)$ over all $q$ in a compact neighborhood of $q_{0}$ and for all $t \leq T$.

Given the vector field $Y$, let $\Omega_{1}=\{Y\}$, and define recursively the set $\Omega_{k}$ to be the collection of vector fields $-\frac{1}{2}\left\langle\bar{B}_{i}: \bar{B}_{k-i}\right\rangle$, for all $B_{i} \in \Omega_{i}$ and $B_{k-i} \in \Omega_{k-i}$. The first few sets are:

$$
\left.\begin{array}{c}
\Omega_{1}=\{Y\}, \quad \Omega_{2}=\left\{-\frac{1}{2}\langle\bar{Y}: \bar{Y}\rangle\right\}, \quad \Omega_{3}=\left\{\frac{1}{4}\langle\bar{Y}: \overline{\langle\bar{Y}: \bar{Y}\rangle}\rangle\right\}  \tag{3.16}\\
\Omega_{4}=\left\{-\frac{1}{8}\langle\bar{Y}: \overline{\langle\bar{Y}}: \overline{\langle\bar{Y}: \bar{Y}\rangle}\rangle\right.
\end{array},-\frac{1}{8}\langle\overline{\langle\bar{Y}: \bar{Y}\rangle}: \overline{\langle\bar{Y}: \bar{Y}\rangle}\rangle\right\} . .
$$

Next, we prove by induction that, for all $k$, the vector field $Y_{k}$ is the sum of $N_{k}$ vector fields belonging to $\Omega_{k}$. The statement is true at $k=1$ with $N_{1}=1$. We assume it true for all $j<k$ and prove it for $k$. Because of the induction assumption, we write $Y_{j}=\sum_{a=1}^{N_{j}} B_{j, a}$, where the $B_{j, a}$ are elements in $\Omega_{j}$. We compute

$$
\begin{aligned}
Y_{k} & =-\frac{1}{2} \sum_{j=1}^{k-1}\left\langle\bar{Y}_{j}: \bar{Y}_{k-j}\right\rangle \\
& =-\frac{1}{2} \sum_{j=1}^{k-1}\left\langle\sum_{a=1}^{N_{j}} \overline{B_{j, a}}: \sum_{b=1}^{N_{k-j}} \overline{B_{k-j, b}}\right\rangle \\
& =\sum_{j=1}^{k-1} \sum_{a=1}^{N_{j}} \sum_{b=1}^{N_{k-j}} \underbrace{-\frac{1}{2}\left\langle\overline{B_{j, a}}: \overline{B_{k-j, b}}\right\rangle}_{\in \Omega_{k}}
\end{aligned}
$$

This concludes the proof by induction and the recursive relation on $N_{k}$ is

$$
\begin{equation*}
N_{1}=1, \quad N_{k}=\sum_{j=1}^{k-1} N_{j} N_{k-j}, \quad k \geq 2 \tag{3.17}
\end{equation*}
$$

As we discuss in Appendix A equation (A.1), the sequence $N_{k}$ can be explicitly computed and bounded as

$$
\begin{equation*}
N_{k}=\frac{1}{k}\binom{2 k-2}{k-1} \leq \frac{2^{2(k-1)}}{k-\frac{1}{2}} \tag{3.18}
\end{equation*}
$$

We now focus our attention on bounding the generic time-varying vector field $(q, s) \mapsto B_{k}(q, s)$ in $\Omega_{k}$. Recall that the symbols $\delta,\|\cdot\|_{\sigma},\left\|\partial^{m} \cdot\right\|_{\sigma}$ introduced in Section 2.7. We claim that there exist sequences of real and integer coefficients $\left\{c_{k}\right.$ : $k \in \mathbb{R}\}$ and $\left\{d_{k}: k \in \mathbb{N}\right\}$, such that

$$
\begin{equation*}
\left\|\partial^{m} B_{k}\right\|_{\sigma^{\prime}} \leq c_{k} \frac{\left(m+d_{k}\right)!}{d_{k}!} \delta^{m+k-1}\|Y\|_{\sigma}^{k} t^{2(k-1)} \tag{3.19}
\end{equation*}
$$

For convenience, we redefine $\delta$ to $\delta=\max \left\{\frac{n}{\sigma-\sigma^{\prime}},\|\Gamma\|_{\sigma}\right\}$, so that $\left\|\partial^{m} \Gamma\right\|_{\sigma^{\prime}} \leq m!\delta^{m+1}$. As discussed in that section, the bound in equation (3.19) is satisfied at $k=1$ for all $m \in \mathbb{N}$, with $c_{1}=1, d_{1}=0$. In what follows we provide a proof by induction on $k \geq 2$.

Any time-varying vector field $B_{k}$ at $k \geq 2$ can be written as $B_{k}=-\frac{1}{2}\left\langle\overline{B_{a}}: \overline{B_{b}}\right\rangle$ for some $1 \leq a, b \leq k-1, a+b=k$ and $B_{a} \in \Omega_{a}, B_{b} \in \Omega_{b}$. Accordingly, we compute:

$$
\begin{aligned}
\| \partial^{m} & \left\langle\overline{B_{a}}: \overline{B_{b}}\right\rangle \|_{\sigma^{\prime}} \\
& =\max _{i, i_{1}, \ldots, i_{m}} \| \frac{\partial^{m}}{\partial q_{i_{1}} \cdots \partial q_{i_{m}}}\left(\frac{\partial \overline{B_{a}^{i}}}{\partial q^{j}} \overline{B_{b}^{j}}+\frac{\partial \overline{B_{b}^{i}}}{\partial q^{j}} \overline{B_{a}^{j}}+\Gamma_{j l}^{i} \overline{\left.\left(\overline{B_{a}^{j}} \overline{B_{b}^{l}}+\overline{B_{b}^{j}} \overline{B_{a}^{l}}\right)\right) \|_{\sigma^{\prime}}} \begin{array}{rl} 
& \leq \max _{i}\left(\left\|\partial^{m}\left(\frac{\partial \overline{B_{a}^{i}}}{\partial q^{j}} \overline{B_{b}^{j}}\right)\right\|_{\sigma^{\prime}}+\left\|\partial^{m}\left(\frac{\partial \overline{B_{b}^{i}}}{\partial q^{j}} \overline{B_{a}^{j}}\right)\right\|_{\sigma^{\prime}}+2\left\|\partial^{m}\left(\Gamma_{j l}^{i} \overline{B_{a}^{j}} \overline{B_{b}^{l}}\right)\right\|_{\sigma^{\prime}}\right) .
\end{array} . . . \begin{array}{l}
\end{array}\right) .
\end{aligned}
$$

Relying on the equality

$$
\frac{d^{m}}{d x^{m}} f(x) g(x)=\sum_{\alpha=0}^{m}\binom{m}{\alpha} \frac{d^{\alpha} f(x)}{d x^{\alpha}} \frac{d^{m-\alpha} g(x)}{d x^{m-\alpha}}
$$

the first term is bounded according to:

$$
\begin{aligned}
&\left\|\partial^{m}\left(\frac{\partial \overline{B_{a}^{i}}}{\partial q^{j}} \overline{B_{b}^{j}}\right)\right\|_{\sigma^{\prime}} \leq n \sum_{\alpha=0}^{m} \frac{m!}{\alpha!(m-\alpha)!}\left\|\partial^{\alpha+1} \overline{B_{a}}\right\|_{\sigma^{\prime}}\left\|\partial^{m-\alpha} \overline{B_{b}}\right\|_{\sigma^{\prime}} \\
& \leq n \sum_{\alpha=0}^{m} \frac{m!}{\alpha!(m-\alpha)!}\left(c_{a} \frac{\left(\alpha+1+d_{a}\right)!}{d_{a}!} \delta^{\alpha+a}\|Y\|_{\sigma}^{a} \frac{t^{2 a-1}}{2 a-1}\right) \\
& \cdot\left(c_{b} \frac{\left(m-\alpha+d_{b}\right)!}{d_{b}!} \delta^{m-\alpha+b-1}\|Y\|_{\sigma}^{b} \frac{t^{2 b-1}}{2 b-1}\right) \\
&= \frac{n c_{a} c_{b}}{(2 a-1)(2 b-1)}\left(\frac{m!}{d_{a}!d_{b}!} \sum_{\alpha=0}^{m} \frac{\left(\alpha+1+d_{a}\right)!\left(m-\alpha+d_{b}\right)!}{\alpha!(m-\alpha)!}\right) \\
& \cdot \delta^{m+a+b-1}\|Y\|_{\sigma}^{a+b} t^{2(a+b-1)}
\end{aligned}
$$

The third term is bounded according to:

$$
\begin{aligned}
&\left\|\partial^{m}\left(\Gamma_{j l}^{i} \overline{B_{a}^{j}} \overline{B_{b}^{l}}\right)\right\|_{\sigma^{\prime}}= n^{2} \sum_{\alpha=0}^{m} \sum_{\beta=0}^{\alpha} \frac{m!\left\|\partial^{m-\alpha} \Gamma\right\|_{\sigma^{\prime}}}{(m-\alpha)!\beta!(\alpha-\beta)!}\left\|\partial^{\beta} \overline{B_{a}}\right\|_{\sigma^{\prime}}\left\|\partial^{\alpha-\beta} \overline{B_{b}}\right\|_{\sigma^{\prime}} \\
&= n^{2} \sum_{\alpha=0}^{m} \sum_{\beta=0}^{\alpha} \frac{m!}{(m-\alpha)!\beta!(\alpha-\beta)!}\left((m-\alpha)!\delta^{m-\alpha+1}\right) \\
& \cdot\left(c_{a} \frac{\left(\beta+d_{a}\right)!}{d_{a}!} \delta^{\beta+a-1}\|Y\|_{\sigma}^{a} \frac{t^{2 a-1}}{2 a-1}\right) \\
& \cdot\left(c_{b} \frac{\left(\alpha-\beta+d_{b}\right)!}{d_{b}!} \delta^{\alpha-\beta+b-1}\|Y\|_{\sigma}^{b} \frac{t^{2 b-1}}{2 b-1}\right) \\
&= \frac{n^{2} c_{a} c_{b}}{(2 a-1)(2 b-1)}\left(\frac{m!}{d_{a}!d_{b}!} \sum_{\alpha=0}^{m} \sum_{\beta=0}^{\alpha} \frac{\left(\beta+d_{a}\right)!\left(\alpha-\beta+d_{b}\right)!}{\beta!(\alpha-\beta)!}\right) \\
& \cdot \delta^{m+a+b-1}\|Y\|_{\sigma}^{a+b} t^{2(a+b-1)} .
\end{aligned}
$$

To simplify notation, let us define

$$
S\left(l, d_{1}, d_{2}\right) \triangleq \sum_{a=0}^{l} \frac{\left(a+d_{1}\right)!\left(l-a+d_{2}\right)!}{a!(l-a)!}
$$

Putting it all together:

$$
\begin{aligned}
\left\|\partial^{m}\left\langle\overline{B_{a}}: \overline{B_{b}}\right\rangle\right\|_{\sigma^{\prime}} \leq & \frac{n c_{a} c_{b}}{(2 a-1)(2 b-1)} \frac{m!}{d_{a}!d_{b}!} \delta^{m+a+b-1}\|Y\|_{\sigma}^{a+b} t^{2(a+b-1)} \\
& \cdot\left(S\left(m, d_{a}+1, d_{b}\right)+S\left(m, d_{a}, d_{b}+1\right)+2 n \sum_{\alpha=0}^{m} S\left(\alpha, d_{a}, d_{b}\right)\right)
\end{aligned}
$$

Equation (A.2) in Appendix A implies that

$$
\begin{equation*}
S\left(l, d_{1}, d_{2}\right)=\frac{d_{1}!d_{2}!\left(l+1+d_{1}+d_{2}\right)!}{l!\left(1+d_{1}+d_{2}\right)!} \tag{3.20}
\end{equation*}
$$

so that we compute

$$
\begin{array}{r}
\frac{m!}{d_{a}!d_{b}!}\left(S\left(m, d_{a}+1, d_{b}\right)+S\left(m, d_{a}, d_{b}+1\right)+2 n \sum_{\alpha=0}^{m} S\left(\alpha, d_{a}, d_{b}\right)\right) \\
=\frac{m!}{d_{a}!d_{b}!}\left(\left(\left(d_{a}+1\right)!d_{b}!+d_{a}!\left(d_{b}+1\right)!\right) \frac{\left(m+2+d_{a}+d_{b}\right)!}{m!\left(2+d_{a}+d_{b}\right)!}\right. \\
\left.\quad+2 n \sum_{\alpha=0}^{m} \frac{d_{a}!d_{b}!\left(\alpha+1+d_{a}+d_{b}\right)!}{\alpha!\left(1+d_{a}+d_{b}\right)!}\right) \\
=\frac{\left(m+2+d_{a}+d_{b}\right)!}{\left(1+d_{a}+d_{b}\right)!}+\frac{2 n m!}{\left(1+d_{a}+d_{b}\right)!} \underbrace{\sum_{\alpha=0}^{m} \frac{\left(\alpha+1+d_{a}+d_{b}\right)!}{\alpha!}},
\end{array}
$$

and again applying equation (3.20) with $\left(l, d_{1}, d_{2}\right)=\left(m, 1+d_{a}+d_{b}, 0\right)$

$$
\begin{aligned}
& =\frac{\left(m+2+d_{a}+d_{b}\right)!}{\left(1+d_{a}+d_{b}\right)!}+\frac{2 n m!}{\left(1+d_{a}+d_{b}\right)!} \frac{\left(1+d_{a}+d_{b}\right)!\left(m+2+d_{a}+d_{b}\right)!}{m!\left(2+d_{a}+d_{b}\right)!} \\
& =\frac{\left(m+2+d_{a}+d_{b}\right)!}{\left(2+d_{a}+d_{b}\right)!}\left(2+2 n+d_{a}+d_{b}\right)
\end{aligned}
$$

Substituting in

$$
\begin{aligned}
& \left\|\partial^{m}\left\langle\overline{B_{a}}: \overline{B_{b}}\right\rangle\right\|_{\sigma^{\prime}} \\
& \quad \leq \frac{n c_{a} c_{b}\left(2+2 n+d_{a}+d_{b}\right)\left(m+2+d_{a}+d_{b}\right)!}{(2 a-1)(2 b-1)\left(2+d_{a}+d_{b}\right)!} \delta^{m+a+b-1}\|Y\|_{\sigma}^{a+b} t^{2(a+b-1)}
\end{aligned}
$$

Next, we express everything back in terms of $k=a+b$ and $B_{k}=-\frac{1}{2}\left\langle\overline{B_{a}}: \overline{B_{b}}\right\rangle$. We have that:
$\left\|\partial^{m} B_{k}\right\|_{\sigma^{\prime}} \leq \max _{a+b=k}\left(\frac{n c_{a} c_{b}\left(2+2 n+d_{a}+d_{b}\right)}{2(2 a-1)(2 b-1)} \frac{\left(m+2+d_{a}+d_{b}\right)!}{\left(2+d_{a}+d_{b}\right)!}\right) \delta^{m+k-1}\|Y\|_{\sigma}^{k} t^{2(k-1)}$
Equation (3.19) is proven by defining sequences $c_{k}$ and $d_{k}$ such that $c_{1}=1, d_{1}=0$ together with

$$
\begin{aligned}
& d_{k} \leq \max _{a+b=k} 2+d_{a}+d_{b} \\
& c_{k} \leq \max _{a+b=k} \frac{n c_{a} c_{b}\left(2+2 n+d_{a}+d_{b}\right)}{2(2 a-1)(2 b-1)}
\end{aligned}
$$

It is immediate to see that $d_{k}=2(k-1)$ satisfies the recursive requirement, so that we require $c_{k}$ to satisfy $c_{1}=1$ together with the requirement

$$
c_{k} \leq \max _{a+b=k} \frac{n(k+n-1) c_{a} c_{b}}{(2 a-1)(2 b-1)}=\max _{a \in\{1, \ldots k-1\}} \frac{n(k+n-1) c_{a} c_{k-a}}{(2 a-1)(2 k-2 a-1)}
$$

Consider the polynomial $p(a)=(2 a-1)(2 k-2 a-1)$ in $a \in[1, k-1]$, it assumes its minimum value $(2 k-3)$ at $a=1$, or equivalently $a=k-1$. Accordingly, a stricter requirement on $c_{k}$ is

$$
c_{k} \leq \max _{a \in\{1, \ldots k-1\}} \frac{n(k+n-1)}{2 k-3} c_{a} c_{k-a}
$$

Since $(k-1) /(2 k-3) \leq 1$ and $n /(2 k-3) \leq n$ for all $k \geq 2$, a conservative selection of $c_{k}$ that satisfies this requirement is provided by the sequence

$$
c_{1}=1, \quad c_{k}=n(1+n) \sum_{a=1}^{k-1} c_{a} c_{k-a}, \quad k \geq 2
$$

Recalling the definition in equation (3.17), one can show that $c_{k}=(n(1+n))^{k-1} N_{k}$.
Finally, we summarize all the analysis in Part II and prove convergence. Evaluating at $m=0$ the bound in equation (3.19), we have

$$
\left\|B_{k}\right\|_{\sigma^{\prime}} \leq(n(1+n))^{k-1} N_{k} \delta^{k-1}\|Y\|_{\sigma}^{k} t^{2(k-1)}
$$

and recalling the bound in equation (3.18), we compute

$$
\begin{aligned}
\left\|Y_{k}\right\|_{\sigma^{\prime}} & \leq N_{k}\left\|B_{k}\right\|_{\sigma^{\prime}} \leq(n(1+n))^{k-1} N_{k}^{2} \delta^{k-1}\|Y\|_{\sigma}^{k} t^{2(k-1)} \\
& \leq \frac{\left(2^{4} n(1+n) \delta\right)^{k-1}}{(k-1 / 2)^{2}}\|Y\|_{\sigma}^{k} t^{2(k-1)}
\end{aligned}
$$

An immediate consequence is that for $\left(2^{4} n(n+1) \delta\right)\|Y\|_{\sigma} T^{2}<1$, the series

$$
Y_{\infty}(q, t) \triangleq \lim _{K \rightarrow \infty} \sum_{k=1}^{K} Y_{k}(q, t)
$$

converges absolutely and uniformly in $t \in[0, T]$ and $q \in B_{\sigma^{\prime}}\left(q_{0}\right)$.
Part III: Here we provide the final limiting argument by collecting various results in Part I, Part II and in Lemma 3.2

We start by studying the behavior as $k \rightarrow \infty$ of the equation (3.14) and of the initial value problem (3.15) from Part I. We shall exploit a variation of a standard result on the continuous dependence of solutions of differential equations with respect to parameter changes, see [16, Chapter I, Section 3]. Uniform convergence of the vector field describing a differential equation, say for example $\sum_{k=1}^{K} Y_{k}$, implies the uniform convergence of the solution to the $K$ th differential equation to the solution of the limiting differential equation. In order to apply this result to the differential equation (3.15) we need to ensure that the vector field on right hand side converges uniformly and absolutely.

Assume that the time length $T$ and input vector field $Y$ satisfy the bound in equation (3.4) inside the theorem statement. Then Lemma 3.2 guarantees that $\gamma([0, T]) \subset B_{\sigma^{\prime}}\left(q_{0}\right)$, and the analysis in Part II guarantees that series $\sum_{k=1}^{\infty} Y_{k}$ converges absolutely and uniformly over $q \in B_{\sigma^{\prime}}\left(q_{0}\right)$. Therefore the series converges uniformly and absolutely along the curve $\gamma$. From equation (3.14) one can deduce that $\gamma(t)=\pi\left(v_{q, k}(t)\right)$, so that the series $\sum_{k=1}^{\infty} Y_{k}$ converges also along $\pi \circ v_{q, k}:[0, T] \mapsto Q$. Accordingly, we can take the limit as $k \rightarrow \infty$ in equation (3.15).

Notice that uniformly in $t \in[0, T]$ and $q \in B_{\sigma^{\prime}}\left(q_{0}\right)$

$$
\lim _{k \rightarrow \infty} Y_{k}(q, t)=0, \quad \text { and } \quad \lim _{k \rightarrow \infty} W_{k}(q, t)=0
$$

and define the time-varying vector field

$$
V_{\infty}=\sum_{k=1}^{\infty} V_{k}=\sum_{k=1}^{\infty} \bar{Y}_{k} .
$$

Taking the limit as $k \rightarrow \infty$ in equations (3.14) and (3.15) one obtains

$$
\dot{\gamma}(t)=v_{q, \infty}(t)+V_{\infty}\left(\pi\left(v_{q, \infty}(t)\right), t\right)
$$

where the curve $v_{q, \infty}:[0, T] \mapsto T Q$ is the solution to

$$
\begin{align*}
\frac{d v_{q, \infty}}{d t} & =\left(Z+\left[V_{\infty}^{\text {lift }}, Z\right]\right)\left(v_{q, \infty}, t\right)  \tag{3.21}\\
v_{q, \infty}(0) & =0_{q_{0}} .
\end{align*}
$$

According to the discussion in Section 2.6, the initial value problem in equation (3.21) can be explicitly integrated. Because $Z \in \mathcal{P}_{1},\left[V_{\infty}^{\text {lift }}, Z\right] \in \mathcal{P}_{0}$, and because of the equality

$$
T \pi \circ\left[V_{\infty}^{\mathrm{lift}}, Z\right]=V_{\infty}
$$

the curve $v_{q, \infty}$ satisfies

$$
v_{q, \infty}(t)=0_{\zeta(t)}, \quad \text { where } \quad \zeta(t)=\Phi_{0, t}^{V_{\infty}}\left(q_{0}\right)
$$

and plugging in

$$
\dot{\gamma}(t)=0_{\zeta(t)}+V_{\infty}\left(\pi\left(0_{\zeta(t)}\right), t\right)=V_{\infty}(\zeta(t), t)
$$

The last two statements imply $\gamma=\zeta$ and are equivalent to statement in equation (3.5) inside the theorem.

Two brief comments are appropriate. First, it is interesting to emphasize an intermediate result proved in Part $I I$ : the $V_{k}$ term in the series is the sum of known number of vector fields belonging to the set $\Omega_{k}$, see the definition preceding equation (3.16). This additional structure might be useful in controllability or motion planning studies. Second, it is unpleasant to remark that while the series expansion is stated in a coordinate-free context, its convergence properties rely on the introduction of a coordinate system.
4. Applications and Extensions. We present a few diverse comments in order to relate the theorem to various earlier works, as well as obtain stronger results under specific additional assumptions on the system.
4.1. The first few order terms and small amplitude forcing. Equation (3.5) is well-defined in the sense that, at fixed $q$, the integration is performed with respect to the time variable. Using the abbreviated notation introduced in equation (2.1),
the first few terms of the sequence $\left\{V_{k}: k \in \mathbb{N}\right\}$ are computed as

$$
\begin{aligned}
& V_{1}=\bar{Y} \\
& V_{2}=-\frac{1}{2} \overline{\langle\bar{Y}: \bar{Y}\rangle} \\
& V_{3}=\frac{1}{2} \overline{\langle\overline{\langle\bar{Y}: \bar{Y}\rangle}: \bar{Y}\rangle} \\
& \left.V_{4}=-\frac{1}{2} \overline{\langle\overline{\langle\overline{\langle\bar{Y}}: \bar{Y}\rangle}: \bar{Y}\rangle}: \bar{Y}\right\rangle \\
& -\frac{1}{8} \overline{\langle\overline{\langle\bar{Y}: \bar{Y}\rangle}: \overline{\langle\bar{Y}: \bar{Y}\rangle}\rangle}
\end{aligned}
$$

so that we can write

$$
\begin{aligned}
\dot{\gamma}(t)= & \bar{Y}(\gamma, t)-\frac{1}{2} \overline{\langle\bar{Y}: \bar{Y}\rangle}(\gamma, t)+\frac{1}{2} \overline{\langle\overline{\langle\bar{Y}: \bar{Y}\rangle}: \bar{Y}\rangle}(\gamma, t) \\
& -\frac{1}{2} \overline{\langle\overline{\langle\overline{\langle\bar{Y}: \bar{Y}\rangle}: \bar{Y}\rangle}: \bar{Y}\rangle}(\gamma, t)-\frac{1}{8} \overline{\langle\overline{\langle\bar{Y}: \bar{Y}\rangle}: \overline{\langle\bar{Y}: \bar{Y}\rangle}\rangle}(\gamma, t)+O\left(\|Y\|^{5} t^{9}\right) .
\end{aligned}
$$

This series converges under the assumption that the product of final time $T$ and input magnitude $\|Y\|_{\sigma}$ be small. Typically, in controllability studies [42] it is the final time that is assumed to be small (the famous acronym STLC stands for small-time local controllability). Within the context of motion planning problems [31, 14], it is instead convenient to study the small input magnitude case. Motivated by the treatment in [14], let $\epsilon$ be a small positive constant, and consider a total acceleration of the form

$$
Y(q, t, \epsilon)=\epsilon X_{1}(q, t)+\epsilon^{2} X_{2}(q, t)+\epsilon^{3} X_{3}(q, t), \quad t \in[0,1]
$$

Accordingly, equation (3.5) is equivalent to

$$
\begin{aligned}
\dot{\gamma}(t)=\epsilon \bar{X}_{1}(\gamma, t)+ & \epsilon^{2}\left(\bar{X}_{2}-\frac{1}{2} \overline{\left\langle\bar{X}_{1}: \bar{X}_{1}\right\rangle}\right)(\gamma, t) \\
& +\epsilon^{3}\left(\bar{X}_{3}-\frac{1}{2} \overline{\left\langle\overline{\left.X_{1}: \bar{X}_{2}\right\rangle}+\right.}+\frac{1}{2} \overline{\left\langle\overline{\left\langle\bar{X}_{1}: \bar{X}_{1}\right\rangle}: \bar{X}_{1}\right\rangle}\right)(\gamma, t)+O\left(\epsilon^{4}\right)
\end{aligned}
$$

This expression generalizes the results presented in Proposition 4.1 in [14]. Note that those results were proven via a perturbation theory argument that is not as general and powerful as the treatment in Theorem 3.3.
4.2. Simple Hamiltonian systems with integrable forces. In this and the following section we analyze systems with more structure both in the affine connections $\nabla$ as well as in the input forces $Y$. Here we consider systems with Lagrangian equal to "kinetic minus potential" and with integrable forces. In the interest of brevity, we refer to the textbooks $[17,35]$ for a detailed presentation and review here only the necessary notation. The affine connection of a simple system is the Levi-Civita connection associated with the kinetic energy matrix $M$, that is, the Christoffel symbols are defined according to the usual relationship

$$
\begin{equation*}
\Gamma_{i j}^{k}=\frac{1}{2} M^{m k}\left(\frac{\partial M_{m j}}{\partial q^{i}}+\frac{\partial M_{m i}}{\partial q^{j}}-\frac{\partial M_{i j}}{\partial q^{m}}\right) \tag{4.1}
\end{equation*}
$$

where $M_{i j}$ and $M^{m k}$ are the components of the matrix representation of $M$ and of its inverse. An integrable time-varying force is written as

$$
\begin{equation*}
Y(q, t)=\operatorname{grad} \varphi(q, t), \quad \text { where } \quad(\operatorname{grad} \varphi)^{i}=M^{i j} \frac{\partial \varphi}{\partial q_{j}} \tag{4.2}
\end{equation*}
$$

and where $\varphi$ is a scalar function on $\mathbb{R}^{n} \times \mathbb{R}$.
One remarkable simplification takes place for a simple system described by a Levi-Civita connection: the set of gradient vector fields is closed under the operation of symmetric product. Let $\varphi_{1}, \varphi_{2}$ be scalar functions on $\mathbb{R}^{n}$ and define a symmetric product between functions according to

$$
\begin{equation*}
\left\langle\varphi_{1}: \varphi_{2}\right\rangle \triangleq \frac{\partial \varphi_{1}}{\partial q} M^{-1} \frac{\partial \varphi_{2}}{\partial q} \tag{4.3}
\end{equation*}
$$

Then the symmetric product of the corresponding gradient vector fields equals the gradient of the symmetric product of the functions. In equations:

$$
\left\langle\operatorname{grad} \varphi_{1}: \operatorname{grad} \varphi_{2}\right\rangle=\operatorname{grad}\left\langle\varphi_{1}: \varphi_{2}\right\rangle
$$

We refer to [10] for the proof. Accordingly, the main theorem can be restated as follows.

Theorem 4.1. Consider the system as described in Problem 3.1. Additionally, let the Christoffel symbols and the input vector field be defined as in equation (4.1) and (4.2). Define recursively the time-varying functions:

$$
\begin{aligned}
& \varphi_{1}(q, t)=\int_{0}^{t} \varphi(q, s) d s \\
& \varphi_{k}(q, t)=-\frac{1}{2} \sum_{j=1}^{k-1} \int_{0}^{t}\left\langle\varphi_{j}(q, s): \varphi_{k-j}(q, s)\right\rangle d s, \quad k \geq 2
\end{aligned}
$$

Then the solution $\gamma:[0, T] \rightarrow Q$ satisfies

$$
\begin{equation*}
\dot{\gamma}(t)=\operatorname{grad} \sum_{k=1}^{+\infty} \varphi_{k}(\gamma(t), t) \tag{4.4}
\end{equation*}
$$

In other words, the flow of a simple Hamiltonian system forced from rest is written as a (time-varying) gradient flow. For completeness, we include a convergence treatment derived from the one in the main theorem.

REMARK 4.2. Given $0<\sigma^{\prime \prime}<\sigma^{\prime}<\sigma$, we assume $M$ and $\varphi$ to be analytic in a neighborhood $B_{\sigma}\left(q_{0}\right)$ of $q_{0}$ and uniformly integrable in $t \in[0, T]$. Two immediate bounds are

$$
\begin{aligned}
\|\operatorname{grad} \varphi\|_{\sigma^{\prime}} & \leq n\left\|M^{-1}\right\|_{\sigma^{\prime}}\left\|\frac{\partial \varphi}{\partial q}\right\|_{\sigma^{\prime}} \\
\|\Gamma\|_{\sigma^{\prime}} & \leq A \triangleq \frac{3 n^{2}}{2\left(\sigma-\sigma^{\prime}\right)}\left\|M^{-1}\right\|_{\sigma^{\prime}}\|M\|_{\sigma}
\end{aligned}
$$

Accordingly, the bounds in the main theorem can be restated (in a more conservative manner) as follows. If

$$
\left\|\frac{\partial \varphi}{\partial q}\right\|_{\sigma^{\prime}} T^{2}<\frac{1}{n\left\|M^{-1}\right\|_{\sigma^{\prime}}} \min \left\{\frac{\sigma^{\prime}-\sigma^{\prime \prime}}{2^{4} n^{2}(n+1)}, \frac{1}{2^{4} n(n+1) A}, \frac{\eta^{2}\left(n^{2} \sigma^{\prime \prime} A\right)}{n^{2} A}\right\}
$$

the series $\sum_{k=1}^{\infty} \varphi_{k}(q, t)$ converges absolutely and uniformly in $t$ and $q$ for all $t \in[0, T]$ and for all $q$ in a neighborhood $B_{\sigma^{\prime \prime}}\left(q_{0}\right)$ of $q_{0}$.
4.3. Invariant systems on Lie groups. In this section we briefly investigate systems with kinetic energy and input forces invariant under a certain group action. These systems have a configuration space $G$ with the structure of an $n$ dimensional matrix Lie group. Systems in this class include satellites, hovercraft, and underwater vehicles.

The equation of motion (2.7) decouples into a kinematic and dynamic equation in the configuration variable $g \in G$ and the body velocity ${ }^{3} v \in \mathbb{R}^{n}$. The kinematic equation can be written ${ }^{4}$ as a matrix differential equation using matrix group notation $\dot{g}=g \widehat{v}$; we refer to [36] for the details. The dynamic equation, sometimes referred to as Euler-Poincarè, is

$$
\begin{equation*}
\dot{v}^{i}+\gamma_{j k}^{i} v^{j} v^{k}=y^{i}(t) \tag{4.5}
\end{equation*}
$$

where the coefficients $\gamma_{j k}^{i}$ are determined by the group and metric structure. The curve $y:[0, T] \mapsto \mathbb{R}^{n}$ denotes the time-varying forcing.

Within this setting, the result in Theorem 3.3 is summarized as follows. The solution to the equation (4.5) with initial condition $v(0)=0$ is $v(t)=\sum_{k=1}^{\infty} v_{k}(t)$, where

$$
\begin{aligned}
& v_{1}(t)=\int_{0}^{t} y(s) d s \\
& v_{k}(t)=-\frac{1}{2} \sum_{j=1}^{k-1} \int_{0}^{t}\left\langle v_{j}(s): v_{k-j}(s)\right\rangle d s, \quad k \geq 2
\end{aligned}
$$

and where the symmetric product between velocity vectors is $\langle x: y\rangle^{i}=-2 \gamma_{j k}^{i} x^{j} y^{k}$. Local convergence for the series expansion can be easily established in this setting.

This result agrees and indeed supersedes the ones presented in [14] obtained via the perturbation method. The relationship of this case to the more general setting studied in Theorem 3.3 is clarified via the notion of invariant connection; see [5, Appendix B] and [40, Section 27, "Variations on a theme by Euler"] for more details.
4.4. Simulations for a three degree of freedom manipulator. In this section we illustrate the approximations derived in Theorem 3.3 by applying them to an example system. We consider a three link planar manipulator. The configuration is described by three angles $\left(\theta_{1}, \theta_{2}, \theta_{3}\right)$. A constant (integrable) force is applied to the first variable. Specifically, we set $\varphi(q, t)=\epsilon \theta_{1}$ and we let the parameter $\epsilon$ vary in the range $10^{-2}$ to 1 . The integration time is $T=1$ seconds. Setting all lengths, masses and moments of inertia to unity, the kinetic energy matrix is:

$$
M=\frac{1}{16}\left[\begin{array}{ccc}
25 & 6 \cos \left(\theta_{1}-\theta_{2}\right) & 2 \cos \left(\theta_{1}-\theta_{3}\right) \\
6 \cos \left(\theta_{1}-\theta_{2}\right) & 21 & 2 \cos \left(\theta_{2}-\theta_{3}\right) \\
2 \cos \left(\theta_{1}-\theta_{3}\right) & 2 \cos \left(\theta_{2}-\theta_{3}\right) & 17
\end{array}\right]
$$

The initial condition is assumed to be $q(0)=(0, \pi / 4,0)$. We investigate the error value $e_{\epsilon, N}=\left\|\gamma(T)-\gamma_{N}(T)\right\|$, where $\gamma_{N}$ is the solution to the $N$ th order truncation: $\dot{\gamma}_{N}(t)=\operatorname{grad} \sum_{k=1}^{N} \varphi_{k}\left(\gamma_{N}, t\right)$. An empirical forecast of the $e_{\epsilon, N}$ is computed as

[^3]Table 4.1
Numerical comparison of various degrees of approximations. The entries in the table are the error values $e_{\epsilon, N}$ that provide a measure of the accuracy of the $N$ th order truncated approximation.

| $\epsilon$ | 1 | .1 | .01 |
| :---: | :---: | :--- | :--- |
| $N=1$ | $5.3 \cdot 10^{-3}$ | $4.6 \cdot 10^{-5}$ | $4.5 \cdot 10^{-7}$ |
| $N=2$ | $2.4 \cdot 10^{-4}$ | $4.2 \cdot 10^{-7}$ | $4.2 \cdot 10^{-10}$ |
| $N=3$ | $1.4 \cdot 10^{-4}$ | $3.0 \cdot 10^{-9}$ | $2.3 \cdot 10^{-13}$ |
| $N=4$ | $5.2 \cdot 10^{-5}$ | $2.4 \cdot 10^{-10}$ | $3.5 \cdot 10^{-15}$ |

follows. Since $T=1$ and $\|Y\|=O(\epsilon)$, there exist two constants $c, d$ such that the $k$ th term in series is bounded by $c(d \epsilon)^{k}$. Summing the neglected contributions from $k=N+1$ to infinity and assuming that $d \epsilon \ll 1$, one can compute $e_{\epsilon, N} \approx c(d \epsilon)^{N+1}$. We summarize the results of the numerical investigation ${ }^{5}$ in Table 4.1. The results are in qualitative agreement with the theoretical forecasts.
5. Conclusions. We have presented a series expansion that describes the evolution of a forced mechanical system. Our result provides a first order description to the solutions of a second order initial value problem. Both the series and the proof method provide insight into the geometry of mechanical control systems. The treatment expands on our previous work [10] on high amplitude high frequency averaging and vibrational stabilization.

Series expansions are the underlying technique for controllability and motion planning. For mechanical systems moving in the low velocity regime, these two problems have been tackled with various degrees of success in [33, 13]. Future research will rely on the contributions in this work to develop more general motion planning algorithms than the ones in [14], and sharper sufficient controllability tests than the ones in [33].

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Appendix A. Some Basic Identities in Combinatorial Analysis. We here present a basic result and derive a useful expression that is needed in the proof of the main theorem. The main reference is the method of generating functions as described in Section 3.4 in [4]. The first identity is explicitly proven in the reference. If $N_{1}=1$ and

$$
N_{k}=\sum_{j=1}^{k-1} N_{j} N_{k-j}, \quad k \geq 2,
$$

then

$$
\begin{equation*}
N_{k}=\frac{1}{k}\binom{2 k-2}{k-1}=\frac{2^{k}}{(4 k-2)} \frac{1 \cdot 3 \cdot 5 \cdots(2 k-1)}{1 \cdot 2 \cdot 3 \cdots k} \leq \frac{4^{k}}{4 k-2} . \tag{A.1}
\end{equation*}
$$

The second equality needed in the proof of Theorem 3.3 is

$$
\begin{equation*}
\sum_{a=0}^{k}\binom{a+d_{1}}{d_{1}}\binom{k-a+d_{2}}{d_{2}}=\binom{k+1+d_{1}+d_{2}}{k} . \tag{A.2}
\end{equation*}
$$

To prove it we use the method of generating functions, see [4]. We claim that, for all real $x$ with $|x|<1$

$$
\begin{equation*}
\sum_{k=0}^{\infty}\left(\sum_{a=0}^{k}\binom{a+d_{1}}{d_{1}}\binom{k-a+d_{2}}{d_{2}}\right) x^{k}=\sum_{k=0}^{\infty}\binom{k+1+d_{1}+d_{2}}{k} x^{k} \tag{A.3}
\end{equation*}
$$

The first step is to notice that

$$
\begin{aligned}
\sum_{k=0}^{\infty}\left(\sum_{a=0}^{k}\binom{a+d_{1}}{d_{1}}\binom{k-a+d_{2}}{d_{2}}\right) x^{k} & =\sum_{k=0}^{\infty} \sum_{m+n=k}\binom{n+d_{1}}{d_{1}}\binom{m+d_{2}}{d_{2}} x^{n+m} \\
& =\left(\sum_{n=0}^{\infty}\binom{n+d_{1}}{d_{1}} x^{n}\right)\left(\sum_{m=0}^{\infty}\binom{m+d_{2}}{d_{2}} x^{m}\right) .
\end{aligned}
$$

Accordingly, we define

$$
\begin{equation*}
f_{a}(x)=\sum_{m=0}^{\infty}\binom{m+a}{a} x^{m} \tag{A.4}
\end{equation*}
$$

and the thesis in equation (A.3) is equivalent to proving that

$$
\begin{equation*}
f_{d_{1}}(x) f_{d_{2}}(x)=f_{d_{1}+d_{2}+1}(x) \tag{A.5}
\end{equation*}
$$

In passing, we also note that the convergence radius of $f$ is $|x|<1$.
The second step is to study the properties of $f$. First of all,

$$
\begin{aligned}
f_{a}(x) & =\sum_{m=0}^{\infty} \frac{(m+a)!}{m!a!} x^{m}=\frac{1}{a!} \sum_{m=0}^{\infty}(m+a) \cdots(m+1) x^{m} \\
& =\frac{1}{a!} \frac{d^{a}}{d x^{a}} \sum_{m=0}^{\infty} x^{m+a}=\frac{1}{a!} \frac{d^{a}}{d x^{a}}\left(x^{a} \sum_{m=0}^{\infty} x^{m}\right)=\frac{1}{a!} \frac{d^{a}}{d x^{a}} \frac{x^{a}}{1-x} .
\end{aligned}
$$

Additionally, it is immediate to see that

$$
f_{0}(x)=\frac{1}{1-x}, \quad x f_{0}(x)=f_{0}(x)-1
$$

and consequently

$$
f_{a}(x)=\frac{1}{a!} \frac{d^{a}}{d x^{a}} \frac{1}{1-x}=\frac{1}{a!} \frac{d^{a}}{d x^{a}} f_{0}(x)
$$

Finally, we prove by induction that

$$
\begin{equation*}
f_{a}(x)=f_{0}(x)^{a+1} \tag{A.6}
\end{equation*}
$$

At $a=0$ the statement is obvious. We assume it true up to $a$ and compute:

$$
\begin{aligned}
f_{a+1}(x) & =\frac{1}{(a+1)!} \frac{d^{a+1}}{d x^{a+1}} \frac{1}{1-x}=\frac{1}{(a+1)!} \frac{d^{a}}{d x^{a}}\left(\frac{1}{1-x}\right)^{2} \\
& =\frac{1}{(a+1)!} \sum_{b=0}^{a}\binom{a}{b}\left(\frac{d^{b}}{d x^{b}} \frac{1}{1-x}\right)\left(\frac{d^{a-b}}{d x^{a-b}} \frac{1}{1-x}\right) \\
& =\frac{a!}{(a+1)!} \sum_{b=0}^{a}\left(\frac{1}{b!} \frac{d^{b}}{d x^{b}} \frac{1}{1-x}\right)\left(\frac{1}{(a-b)!} \frac{d^{a-b}}{d x^{a-b}} \frac{1}{1-x}\right) \\
& =\frac{1}{a+1} \sum_{b=0}^{a}\left(\frac{1}{1-x}\right)^{b+1}\left(\frac{1}{1-x}\right)^{a-b+1}=\left(\frac{1}{1-x}\right)^{a+2}
\end{aligned}
$$

This concludes the proof of equation (A.6), which immediately implies equation (A.5) and the main thesis in equation (A.3).


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[^1]:    ${ }^{1}$ We here refer to the $\Gamma_{i j}^{k}$ functions as Christoffel symbols even without requiring $\nabla$ to be a Levi-Civita connection.

[^2]:    ${ }^{2}$ Geometric homogeneity corresponds to the existence of an (infinitesimal) symmetry in the equations of motion. For control systems described by an affine connection the symmetry is invariance under affine time-scaling transformations.

[^3]:    ${ }^{3}$ More precisely, the body velocity $v$ lives in the Lie algebra of the group $G$.
    ${ }^{4}$ Alternatively, the kinematic equation can be written in a system of local coordinates $q$ (e.g., Euler angles in the case of rotation matrices) as $\dot{q}=J(q) v$, where $J(q)$ is an appropriate Jacobian matrix.

[^4]:    ${ }^{5}$ The numerical integration is performed inside the Mathematica environment, specifying 16 digits of accuracy and 32 digits of working precision.

