Fortunati, A., \& Wiggins, S. R. (2016). Integrability and strong normal forms for non-autonomous systems in a neighbourhood of an equilibrium. Journal of Mathematical Physics, 57(9), [092703]. DOI: 10.1063/1.4962802

Early version, also known as pre-print

Link to published version (if available):
10.1063/1.4962802

Link to publication record in Explore Bristol Research
PDF-document

This is the submitted manuscript. The final published version (version of record) is available online via AIP at http://scitation.aip.org/content/aip/journal/jmp/57/9/10.1063/1.4962802. Please refer to any applicable terms of use of the publisher.

## University of Bristol - Explore Bristol Research

## General rights

This document is made available in accordance with publisher policies. Please cite only the published version using the reference above. Full terms of use are available:
http://www.bristol.ac.uk/pure/about/ebr-terms.html

# Integrability and strong normal forms for non-autonomous systems in a neighbourhood of an equilibrium* 

Alessandro Fortunati ${ }^{(a)}$<br>Stephen Wiggins ${ }^{(b)}$<br>School of Mathematics, University of Bristol, Bristol BS8 1TW, United Kingdom


#### Abstract

The paper deals with the problem of existence of a convergent "strong" normal form in the neighbourhood of an equilibrium, for a finite dimensional system of differential equations with analytic and time-dependent non-linear term. The problem can be solved either under some non-resonance hypotheses on the spectrum of the linear part or if the non-linear term is assumed to be (slowly) decaying in time. This paper "completes" a pioneering work of Pustil'nikov in which, despite under weaker non-resonance hypotheses, the nonlinearity is required to be asymptotically autonomous. The result is obtained as a consequence of the existence of a strong normal form for a suitable class of real-analytic Hamiltonians with non-autonomous perturbations.


Keywords: Non-autonomous systems, normal forms, aperiodic time dependence.
2010 MSC: Primary: 37J40. Secondary: 37B55, 37J25.

## 1 Preliminaries and main result

### 1.1 Introduction

The study of the dynamics of a system of ordinary differential equations (ODEs) in a neighbourhood of an equilibrium, boasts nowadays a rich and well established theory. Its foundation goes even back to the late XIX century to the contribution of Poincaré [Poi79] and Lyapunov [Lya92]. Given an analytic vector field, the possibility to write the motions of the associated system in the vicinity of an equilibrium as a convergent power series, is deeply related to some non-resonance conditions on the eigenvalues of the linear part.
The results have been afterwards extended in the studies of Siegel started in [Sie42]. The problem of the reducibility of a given system to a linear form via an analytic transformation, it is shown to be solvable in [Sie52] for a full measure set of eigenvalues .
In the case of Hamiltonian structure, investigated later in [Sie54], the problem can be naturally interpreted in terms of the existence of a (convergent) canonical transformation of variables,

[^0]casting a Hamiltonian of the form "quadratic" + "perturbation" into a suitable ${ }^{1}$ normal form, in some neighbourhood of the examined equilibrium. Based on this approach, the paper [Gio] provides a generalisation of the results by Lyapunov, removing the hypothesis of purely imaginary eigenvalues.
In any case, we remark that, as a common feature of this class of problems, without any assumption on the eigenvalues, the program of casting the Hamiltonian at hand into a normal form, at least in general, fails. In fact, it is immediate to recognize how the linear combinations of eigenvalues occurring in the normalization scheme could produce some "small divisor effects". Knowingly, this phenomenon can either obstruct the formal resolvability of the homological equations produced during the normalization or jeopardize the convergence of the series.
We recall that, for instance, the described problem of well-posedness of the homological equation is overcome by Moser in [Mos56], in the case of "one and a half 2 " degrees of freedom Hamiltonian $H(p, q, t)$ close to a hyperbolic equilibrium located at $p=q=0$. The strategy consists of keeping terms of the form $(p q)^{k}, k \geq 2$, in the normal form. In this way the canonical equations are still integrable ( $x:=p q$ is a prime integral) but this allows to avoid the division by zero in the homological equation which would have been carried by those terms. This analysis plays a fundamental role in the context of instability phenomena in Hamiltonian systems with several degrees of freedom (Arnold's diffusion), in order to describe the flow in the neighbourhood of partially hyperbolic tori of a priori unstable systems, see [CG94].
The pioneering work by Pustil'nikov [Pus74], aims to extend the results of the paper [Sie52], by introducing a time dependence in the non-linear part of the vector field (not necessarily Hamiltonian). As it is natural, the choice of a suitable class of time-dependent perturbations and its treatment is a further difficulty to the phenomenon of the "resonances". In [Pus74], under the non-resonance condition already assumed in [Sie52] for the autonomous case, it is required that the perturbation is asymptotic to a time-independent, analytic function. However, no restrictions are imposed on the "type" of the time dependence, more specifically, it has to be neither periodic nor quasi-periodic. This case is also known as aperiodic time dependence.
After [Pus74], the interest in a general dependence on time has been renewed in [GZ92] then followed by [Bou13], [FW14] and subsequent papers. Basically, all of them deal with the Hamiltonian case (see [FW15a] for the case of Poisson systems). The paper [FW15b] extends the above described result by Moser to the case of a perturbation aperiodically dependent on time.
As a matter of fact, the Hamiltonian structure is not a real obstruction for the use of the tools apt to treat the Hamiltonian case. In fact, given a system of ODEs, it can be always interpreted as ("a half" of the) canonical equations of a suitable Hamiltonian system, of larger dimension, see e.g. [Ber09]. The strategy of this paper is to derive the integrability of the system of ODEs at hand, see (7), as a particular case of the existence of a normal form for a real-analytic Hamiltonian with aperiodic perturbation, see (1), by using the tools introduced in [FW15b] for the one degrees of freedom case.
The possibility to cast the Hamiltonian (1) into a normal form is shown to be possible in the two cases described in Theorem 1.1. In the second case, we deal with perturbations linear in the $y$ variables, in the presence of some non-resonance assumption on the eigenvalues. This case is directly related to the Hamiltonian formulation of a system of ODEs (due to the linearity in $y)$. It is immediate to notice that, with respect to [Pus74, (0.3)], the condition (4) on the eigenvalues is clearly more restrictive. Nevertheless, the hypothesis of asymptotic time-independence assumed in [Pus74] is weakened to the simple boundedness.

[^1]On the other hand, the first case, has a more general character: if the perturbation decays ${ }^{3}$ in time, either the described assumption on the form of $f$ or on the eigenvalues turn out to be unnecessary. Basically, the presence of resonance phenomena is no longer an obstruction for the existence of the normal form, see also [FW15c].
The paper, based on the Lie series formalism developed by A. Giorgilli et al., can be regarded, at the same time, as a non-autonomous version of [Gio].

### 1.2 Setting

Let us consider the following Hamiltonian

$$
\begin{equation*}
H(x, y, \eta, t)=h(x, y, \eta)+f(x, y, t), \quad h(x, y, \eta):=\eta+\sum_{l=1}^{n} \lambda_{l} x_{l} y_{l}, \tag{1}
\end{equation*}
$$

where $(x, y, \eta) \in \mathcal{D}:=[-r, r]^{n} \times[-r, r]^{n} \times \mathbb{R}$, with $n \geq 1$ and $r>0, \lambda_{l} \in \mathbb{C}$ and $t \in \mathbb{R}^{+}:=[0,+\infty)$. The assumptions on $f$ will be discussed below. The system (1) is nothing but the "autonomous equivalent" of $\mathcal{H}(x, y, t)=\sum_{l=1}^{n} \lambda_{l} x_{l} y_{l}+f(y, x, t)$, once $\eta$ has been defined as the conjugate variable to $t$.
The standard use of the analytic tools requires the complexification of the domain $\mathcal{D}$ as follows. Given $R \in(0,1 / 2]$ set $\mathcal{D}_{R}:=\mathcal{Q}_{R} \times \mathcal{S}_{R}$, where

$$
\mathcal{Q}_{R}:=\left\{(x, y) \in \mathbb{C}^{2 n}:|x|,|y| \leq R\right\}, \quad \mathcal{S}_{R}:=\{\eta \in \mathbb{C}:|\Im \eta| \leq R\},
$$

It will be required that, for all $t \in \mathbb{R}^{+}, f$ belongs to the space of real-analytic functions on $\dot{\mathcal{Q}}_{R}$ and continuous on the boundary, which we denote with $\mathfrak{C}\left(\mathcal{Q}_{R}\right)$. In such a way $H \in \mathfrak{C}\left(\mathcal{D}_{R}\right)$.
In particular, the space of all the $G \in \mathfrak{C}\left(\mathcal{Q}_{R}\right)$ is endowed with the Taylor norm

$$
\begin{equation*}
\|G(x, y, t)\|_{R}:=\sum_{\alpha, \beta \in \mathbb{N}^{n}}\left|g_{\alpha, \beta}(t)\right| R^{|\alpha+\beta|} \tag{2}
\end{equation*}
$$

where $G(x, y, t)=: \sum_{\alpha, \beta \in \mathbb{N}^{n}} g_{\alpha, \beta}(t) x^{\alpha} y^{\beta}$ and $^{4}|\alpha|:=\sum_{l=1}^{n} \alpha_{l}$. We recall the standard result for which, if $G \in \mathfrak{C}\left(\mathcal{Q}_{R}\right)$ for all $t \in \mathbb{R}^{+}$, then $\left|g_{\alpha, \beta}(t)\right| \leq|G|_{R} R^{-|\alpha+\beta|}$, where $|G|_{R}:=\sup _{(x, y) \in \mathcal{Q}_{R}}|G|$. In particular, $\|G\|_{R^{\prime}}<+\infty$ for all $R^{\prime}<R$.
Throughout this paper we shall deal with perturbations satisfying the following conditions:

1. $f$ is "at least" quadratic in $x$ and "at least" linear in $y$ : a property that we will denote with $(Q x L y)$, i.e. $f_{\alpha, \beta}(t)=0$ for all $t \in \mathbb{R}^{+}$and for all $(\alpha, \beta) \in \mathbb{N}^{2 n} \backslash \Gamma$, where $\Gamma:=\{(\alpha, \beta) \in$ $\left.\mathbb{N}^{2 n}:|\alpha| \geq 2,|\beta| \geq 1\right\}$,
2. there exist $M_{f} \in[1,+\infty)$ and $a \in[0,1)$ such that ${ }^{5}$, for all $(x, y, t) \in \mathcal{Q}_{R} \times \mathbb{R}^{+}$,

$$
\begin{equation*}
\|f(x, y, t)\|_{R} \leq M_{f} e^{-a t} \tag{3}
\end{equation*}
$$

[^2]
### 1.3 Main result

In the described setting, the main result can be stated as follows
Theorem 1.1. Suppose that one of the following conditions are satisfied:
I. Time decay: $a>0$.
II. Linearity in $y+$ non-resonance: $a=0$ and the perturbation is linear in $y$, denoted by (Ly), i.e. of the form $f(x, y, t)=y \cdot g(x, t)$. In addition, the vector $\Lambda:=\left(\lambda_{1}, \ldots, \lambda_{n}\right)$, satisfies the non-resonance condition

$$
\begin{equation*}
\max _{l=1, \ldots, n}\left(\left|\Re \mathcal{U}\left(\alpha, e_{l}, \Lambda\right)\right|^{-1}\right) \leq \gamma|\alpha|^{\tau}, \quad \forall \alpha \in \mathbb{N}^{n} \tag{4}
\end{equation*}
$$

where $\mathcal{U}(\alpha, \beta, \Lambda):=(\alpha-\beta) \cdot \Lambda$, for some $\gamma>0$ and $\tau \geq n$. $e_{l}$ stands for the $l-t h$ vector of the canonical basis of $\mathbb{R}^{n}$.
Then it is possible to determine $R_{*}, R_{0}$ with $0<R_{*}<R_{0} \leq R^{16}$ and a family of canonical transformations $(x, y, \eta)=\mathcal{M}\left(x^{(\infty)}, y^{(\infty)}, \eta^{(\infty)}\right), \mathcal{M}: \mathcal{D}_{R_{*}} \rightarrow \mathcal{D}_{R_{0}}$, analytic on $\mathcal{D}_{R_{*}}$ for all $t \in \mathbb{R}^{+}$, casting the Hamiltonian (1) into the strong normal form

$$
\begin{equation*}
H^{(\infty)}\left(x^{(\infty)}, y^{(\infty)}, \eta^{(\infty)}\right)=h\left(x^{(\infty)}, y^{(\infty)}, \eta^{(\infty)}\right) \tag{5}
\end{equation*}
$$

Remark 1.2. It is immediate to recognize the similarity between (4) and the standard Diophantine condition. Clearly, all the vectors $\Lambda$ whose real part is a Diophantine vector, satisfy condition (4), no matter what the imaginary part is. Hence the set of vectors satisfying (4) is, $a$ fortiori, a full-measure set.
As anticipated in the introduction, we stress that condition (4) is stronger than the non-resonance condition imposed in [Pus74] and it is not satisfied in the case of purely imaginary $\Lambda$.

Remark 1.3. As usually done in the Lie series method, see e.g. [Gio03], the transformation $\mathcal{M}$ will be constructed as the limit (defined, at the moment, only at a formal level)

$$
\begin{equation*}
\mathcal{M}:=\lim _{j \rightarrow \infty} \mathcal{M}^{(j)} \circ \mathcal{M}^{(j-1)} \circ \ldots \circ \mathcal{M}^{(0)} \tag{6}
\end{equation*}
$$

where $\mathcal{M}^{(j)}:=\exp \left(\mathcal{L}_{\chi^{(j)}}\right) \equiv \operatorname{Id}+\sum_{s \geq 1}(s!)^{-1} \mathcal{L}_{\chi^{(j)}}^{s}$ and $\mathcal{L}_{\chi^{(j)}}:=\left\{\cdot, \chi^{(j)}\right\}$. The generating sequence $\left\{\chi^{(j)}\right\}_{j \in \mathbb{N}}$, where $\chi^{(j)}=\chi^{(j)}(x, y, t)$, see [GZ92], is meant to be determined.
We will show (see the proof of Lemma 3.3) that in the case of a perturbation which is (Ly), it is possible to show that $\chi^{(j)}(x, y, t)$ is $(L y)$ as well, for all $j \in \mathbb{N}$. In such a case, it is easy to check by induction that $x^{(j)}=\mathcal{M}^{(j)} x^{(j+1)}$ does not depend on the variable $y$, for all $j$. Hence the composition $x \equiv x^{(0)}=\mathcal{M} x^{(\infty)}=: \mathcal{M}_{x}\left(x^{(\infty)}, t\right)$ does not depend on $y^{(\infty)}$ i.e. is an analytic map $\mathcal{M}_{x}: \tilde{\mathcal{Q}}_{R_{*}} \rightarrow \tilde{\mathcal{Q}}_{R_{0}}$ parametrised by $t$, where $\tilde{\mathcal{Q}}_{R}:=\left\{x \in \mathbb{C}^{n}:|x| \leq R\right\}$. This will play a key role in the next section.

### 1.4 The corollary

Let us consider the following non-linear system

$$
\begin{equation*}
\dot{v}=A v+g(v, t), \tag{7}
\end{equation*}
$$

where $v \in \mathbb{R}^{n}, A$ is a $n \times n$ matrix with real entries and the function $g$ is such that $\partial_{v}^{\nu} g(0, t) \equiv 0$ for all $\nu \in \mathbb{N}^{n}$ such that $|\nu| \leq 1$ i.e. $g$ is at least quadratic in $v$. We restrict ourselves to the
class of diagonalizable $A$ with non-purely imaginary eigenvalues $\lambda_{l}$. In the obvious system of coordinates denoted with $x$, the system (7) easily reads as

$$
\begin{equation*}
\dot{x}_{l}=\lambda_{l} x_{l}+\tilde{g}_{l}(x, t), \quad l=1, \ldots, n . \tag{8}
\end{equation*}
$$

In this framework one can state the next
Corollary 1.4. Suppose that $f(x, y, t):=y \cdot \tilde{g}(x, t)$ and $\Lambda$ is such that the conditions described in II of Theorem 1.1 are satisfied. Then the system (8) is integrable in a suitable neighbourhood of the origin.
The same result holds, in particular, without any non-resonance condition on $\Lambda$, provided that $\tilde{g}(x, t)$ is such that (3) is satisfied with $a>0$.

Proof. The key remark, see e.g. [Ber09], is that (8) can be interpreted as a set of canonical equations of the Hamiltonian system with Hamiltonian $\mathcal{K}:=\eta+\sum_{l=1}^{n} y_{l}\left(\Lambda_{l} x_{l}+\tilde{g}_{l}(x, t)\right)$, i.e. (1) with $f(x, y, t)$ defined in the statement. Hence, by Theorem 1.1, there exists a suitable neighbourhood of the origin endowed with a set of coordinates $\left(x^{(\infty)}, y^{(\infty)}, \eta^{(\infty)}\right)$, such that $\mathcal{K}$ is cast into the (integrable) strong form $\mathcal{K}^{(\infty)}=\eta^{(\infty)}+\sum_{l=1}^{n} \lambda_{l} y_{l}^{(\infty)} x_{l}^{(\infty)}$. Furthermore, as noticed in Remark 1.3, $\mathcal{M}_{x}$ is an analytic map between $x$ and $x^{(\infty)}$. Hence $x(t)=\mathcal{M}_{x}\left(x^{(\infty)}(0) \exp (\mathcal{A} t), t\right)$, with $\mathcal{A}:=\operatorname{diag}\left(\lambda_{1}, \ldots, \lambda_{n}\right)$, gives the explicit solution of (8).

## 2 Some preliminary results

### 2.1 Two elementary inequalities

Proposition 2.1. For all $\mathcal{R} \leq e^{-4}$ and all $\delta \leq 1 / 2$ the following inequalities hold

$$
\begin{equation*}
\sum_{\substack{\nu \in \mathbb{N}^{m} \\|\nu| \geq N}} \mathcal{R}^{|\nu|} \leq 2 m e^{3 m-3} \mathcal{R}^{\frac{3 N}{4}}, \quad \sum_{\nu \in \mathbb{N}^{m}}|\nu|^{\mu}(1-\delta)^{|\nu|} \leq \mathcal{C}(m, \mu) \delta^{-m-\mu-1} \tag{9}
\end{equation*}
$$

where $m \geq 2, \mu \geq 0$ and $\mathcal{C}(m, \mu):=e^{4 m+\mu-1}(m+\mu)^{(m+\mu)} /(m-1)!$.
Proof. See Appendix.

### 2.2 A result on the homological equation

Proposition 2.2. Consider the following equation

$$
\begin{equation*}
\mathcal{L}_{\chi^{(j)}} h+f^{(j)}=0 \tag{10}
\end{equation*}
$$

where $h$ has been defined in (1) and $f^{(j)}=f^{(j)}(x, y, t)=\sum_{(\alpha, \beta) \in \Gamma} f_{\alpha, \beta}^{(j)}(t) x^{\alpha} y^{\beta}$ satisfies $\left\|f^{(j)}\right\|_{\tilde{R}} \leq$ $M_{j} \exp (-a t)$ for some $a \in[0,1)$. The following statements hold for all $\delta \in(0,1 / 2]$ :

1. If $a>0$, there exists $C_{1}=C_{1}(n, \Lambda)>0$ such that

$$
\begin{equation*}
\left\|\chi^{(j)}\right\|_{(1-\delta) \tilde{R}},\left\|\partial_{t} \chi^{(j)}\right\|_{(1-\delta) \tilde{R}} \leq C_{1} M_{j} a^{-1} \delta^{-2(n+1)} \tag{11}
\end{equation*}
$$

2. If $a=0, f^{(j)}$ is of the form $f^{(j)}=y \cdot g^{(j)}(x, t)$ and $\Lambda$ satisfies (4), there exists $C_{2}=$ $C_{2}(n, \Lambda, \tau, \gamma)>0$ such that

$$
\begin{equation*}
\left\|\chi^{(j)}\right\|_{(1-\delta) \tilde{R}},\left\|\partial_{t} \chi^{(j)}\right\|_{(1-\delta) \tilde{R}} \leq C_{2} M_{j} \delta^{-(n+\tau+2)} \tag{12}
\end{equation*}
$$

Proof. First of all note that $\mathcal{L}_{\chi^{(j)}} h=\partial_{t} \chi^{(j)}+\sum_{l=1}^{n} \lambda_{l}\left(x_{l} \partial_{x_{l}}-y_{l} \partial_{y_{l}}\right) \chi^{(j)}$. By expanding the generating function as $\chi^{(j)}(x, y, t)=\sum_{(\alpha, \beta) \in \mathbb{N}^{2 n}} c_{\alpha, \beta}^{(j)}(t) x^{\alpha} y^{\beta}$, equation (10) reads, in terms of Taylor coefficients, as

$$
\begin{equation*}
\dot{c}_{\alpha, \beta}^{(j)}(t)+\mathcal{U}(\alpha, \beta, \Lambda) c_{\alpha, \beta}^{(j)}=f_{\alpha, \beta}^{(j)}(t) . \tag{13}
\end{equation*}
$$

The solution of (13) is easily written, for all $(\alpha, \beta) \in \Gamma$, as

$$
\begin{equation*}
c_{\alpha, \beta}^{(j)}(t)=e^{-\mathcal{U}(\alpha, \beta, \Lambda) t}\left[c_{\alpha, \beta}^{(j)}(0)+\int_{0}^{t} e^{\mathcal{U}(\alpha, \beta, \Lambda) s} f_{\alpha, \beta}^{(j)}(s) d s\right], \tag{14}
\end{equation*}
$$

while trivially $c_{\alpha, \beta}^{(j)}(t) \equiv 0$ for all $(\alpha, \beta) \in \mathbb{N}^{2 n} \backslash \Gamma$.
Now denote $\mathcal{U}_{R}+i \mathcal{U}_{I}:=\mathcal{U}(\alpha, \beta, \Lambda)$ with $\mathcal{U}_{I, R} \in \mathbb{R}$ and recall that, by hypothesis, $\left|f_{\alpha, \beta}^{(j)}(t)\right| \leq$ $M_{j} \tilde{R}^{-|\alpha+\beta|} e^{-a t}$.
Case $a>0$. For all $(\alpha, \beta) \in \Gamma$ such that $\mathcal{U}_{R} \geq 0$ we choose $c_{\alpha, \beta}^{(j)}(0)=0$ then we have

$$
\left|c_{\alpha, \beta}^{(j)}\right| \leq e^{-\mathcal{U}_{R} t} \int_{0}^{t} e^{\mathcal{U}_{R} s}\left|f_{\alpha, \beta}^{(j)}(s)\right| d s \leq M_{j} \tilde{R}^{-|\alpha+\beta|} \int_{0}^{t} e^{-a s} d s \leq M_{j} \tilde{R}^{-|\alpha+\beta|} a^{-1}
$$

Otherwise, for those $\alpha$ and $\beta$ such that $\mathcal{U}_{R}<0$, redefine $\mathcal{U}_{R}:=-\mathcal{U}_{R}$ with $\mathcal{U}_{R}>0$ and choose $c_{\alpha, \beta}^{(j)}(0):=-\int_{\mathbb{R}^{+}} \exp (\mathcal{U}(\alpha, \beta, \Lambda) s) f_{\alpha, \beta}^{(j)}(s) d s$. Note that $\left|c_{\alpha, \beta}^{(j)}(0)\right|<+\infty$. In this case we have $\left|c_{\alpha, \beta}^{(j)}\right| \leq \exp \left(\mathcal{U}_{R} t\right) \int_{t}^{\infty} \exp \left(-\mathcal{U}_{R} s\right)\left|f_{\alpha, \beta}^{(j)}(s)\right| d s \leq M_{j} \tilde{R}^{-|\alpha+\beta|} a^{-1}$. Hence $\left|c_{\alpha, \beta}^{(j)}\right| \leq M_{j} \tilde{R}^{-|\alpha+\beta|} a^{-1}$ for all $(\alpha, \beta) \in \Gamma$. By recalling (2) one gets $\left\|\chi^{(j)}\right\|_{(1-\delta) \tilde{R}} \leq M_{j} a^{-1} \sum_{(\alpha, \beta) \in \mathbb{N}^{2 n}}(1-\delta)^{|\alpha+\beta|}$. The use of the second of (9) with $\nu:=(\alpha, \beta)$, yields the first part of (11) with $C_{1}$ set for the moment to $\hat{C}_{1}:=\mathcal{C}(2 n, 0)$.
Directly from (13) we get $\left|\dot{c}_{\alpha, \beta}^{(j)}\right| \leq|\alpha+\beta||\Lambda|\left|c_{\alpha, \beta}^{(j)}\right|+\left|f_{\alpha, \beta}^{(j)}\right| \leq a^{-1} M_{j}(1+|\Lambda|)|\alpha+\beta| \tilde{R}^{-|\alpha+\beta|}$. By (9) with $\mu=1$ we get the second of part of (11). The constant is chosen as $C_{1}:=(1+|\Lambda|) \mathcal{C}(2 n, 1)>$ $\hat{C}_{1}$.
Case $a=0$. In such case, the homological equation reads as

$$
\begin{equation*}
\dot{c}_{\alpha, l}^{(j)}(t)+\mathcal{U}\left(\alpha, e_{l}, \Lambda\right) c_{\alpha, l}^{(j)}=f_{\alpha, l}^{(j)}(t) \tag{15}
\end{equation*}
$$

where $f_{\alpha, l}^{(j)}:=\left.f_{\alpha, \beta}^{(j)}\right|_{\beta=e_{l}}$ (the same notation for $\left.c_{\alpha, l}^{(j)}\right)$, for all $\alpha \in \mathbb{N}^{n}$ such that $|\alpha| \geq 2$ and for all $l=1, \ldots, n$. By hypothesis (4), $\mathcal{U}_{R} \neq 0$. Similarly to the case $a>0$, if $\mathcal{U}_{R}>0$ we set $c_{\alpha, l}^{(j)}(0)=0$, otherwhise, $c_{\alpha, l}^{(j)}(0):=-\int_{\mathbb{R}^{+}} \exp \left(\mathcal{U}\left(\alpha, e_{l}, \Lambda\right) s\right) f_{\alpha, l}^{(j)}(s) d s$. Proceeding as before, one obtains, by using (4),

$$
\left|c_{\alpha, l}^{(j)}(t)\right| \leq M_{j} \mathcal{U}_{R}^{-1} \tilde{R}^{-|\alpha|-1} \leq \gamma M_{j}|\alpha|^{\tau} \tilde{R}^{-|\alpha|-1}
$$

This implies $\left\|\chi^{(j)}\right\|_{(1-\delta) \tilde{R}} \leq n \gamma M_{j} \sum_{\alpha \in \mathbb{N}^{n}}|\alpha|^{\tau}(1-\delta)^{|\alpha|}$ which is, by (9), the first part of (12) with $\hat{C}_{2}=n \gamma \mathcal{C}(n, \tau)$. On the other hand, from the homological equation, we get $\left|\dot{c}_{\alpha, l}^{(j)}(t)\right| \leq$ $M_{j}|\alpha|^{\tau+1}(1+\gamma|\Lambda|) \tilde{R}^{-|\alpha|-1}$. Similarly, the latter yields the second part of (12) with $C_{2}:=$ $\max \left\{n(1+\gamma|\Lambda|) \mathcal{C}(n, \tau+1), \hat{C}_{2}\right\}$.

### 2.3 A bound on the Lie operator

Proposition 2.3. Let $F, G$ be two functions such that $\|F\|_{(1-\tilde{d}) \tilde{R}},\|G\|_{(1-\tilde{d}) \tilde{R}}<+\infty$ for some $\tilde{d} \in(0,1 / 4]$ and $\tilde{R}>0$. Then for all $s \in \mathbb{N}$ the following bound holds

$$
\begin{equation*}
\left\|\mathcal{L}_{G}^{s} F\right\|_{(1-2 \tilde{d}) \tilde{R}} \leq e^{-2} s!\left[e^{2}(\tilde{R} \tilde{d})^{-2}\|G\|_{(1-\tilde{d}) \tilde{R}}\right]^{s}\|F\|_{(1-\tilde{d}) \tilde{R}} \tag{16}
\end{equation*}
$$

Proof. Straightforward from [Gio, Sec 3.2] and [Gio03, Lemma 4.2].

## 3 Proof of the main result: convergence of the normal form

### 3.1 Preparation of the domains

Taking into account the domain restriction imposed by Proposition 2.3, the canonical transformations will be constructed of the form $\mathcal{M}_{j}: \mathcal{D}_{R_{j+1}} \rightarrow \mathcal{D}_{R_{j}} \ni\left(x^{(j)}, y^{(j)}, \eta^{(j)}\right)$ (understood $\left.\left(x^{(0)}, y^{(0)}, \eta^{(0)}\right) \equiv(x, y, \eta)\right)$, where $\left\{\mathcal{D}_{R_{j}}\right\}_{j \in \mathbb{N}}$ is a suitable sequence of nested domains. We will also provide another sequence $\left\{\epsilon_{j}\right\}$ which will be used to control the size of the remainder.
Lemma 3.1. Let us consider the following sequences

$$
\begin{equation*}
\epsilon_{j+1}=K a^{-1} d_{j}^{-\sigma} \epsilon_{j}^{2}, \quad R_{j+1}:=\left(1-2 d_{j}\right) R_{j}, \tag{17}
\end{equation*}
$$

with $\epsilon_{j}, R_{j}<1, d_{j} \leq 1 / 4$ and where $\epsilon_{0}, R_{0}, a, K, \sigma>0$ are given. If

$$
\begin{equation*}
\epsilon_{0} \leq \epsilon_{a}:=a(2 \pi)^{-\sigma} K^{-1}, \tag{18}
\end{equation*}
$$

then it is possible to construct $\left\{d_{j}\right\}_{j \in \mathbb{N}}$ in such a way $R_{j} \geq R_{*}:=R_{0} / 2$ and $\epsilon_{j} \rightarrow 0$ monotonically as $j \rightarrow \infty$.
Remark 3.2. The property $R_{*}>0$ is crucial, as $R_{*}$ is the lower bound for the analyticity radius of the normalised Hamiltonian.

Proof. Straightforward from [FW15c, Lemma 4.4]. We recall that a suitable choice is $\epsilon_{j}=$ $\epsilon_{0}(j+1)^{-\sigma}$, then, by (17), $d_{j}=\left(\epsilon_{0} K a^{-1}\right)^{(1 / \sigma)}(j+2)^{2} /(j+1)^{4}$. From the latter, one has

$$
\begin{equation*}
\sum_{j \geq 0} d_{j} \leq 1 / 6, \tag{19}
\end{equation*}
$$

provided that condition (18) is satisfied.

### 3.2 Iterative lemma

Let us define for all $j \geq 0, H^{(j+1)}:=\mathcal{M}_{j} H^{(j)}$ with $H^{(0)}:=H$.
Lemma 3.3. Under the same hypotheses of Theorem 1.1 and under the condition (18) it is possible to find a $R_{0}$ and a sequence $\left\{\chi^{(j)}\right\}_{j \in \mathbb{N}}$ such that $H^{(j)}(x, y, \eta, t)=h(x, y, \eta)+f^{(j)}(x, y, t)$ with $f^{(j)}$ (QxLy) and such that $\left\|f^{(j)}\right\|_{R_{j}} \leq \epsilon_{j} e^{-a t}$ for all $j$, where $\epsilon_{j}, R_{j}$ are given by (17).

The stated result exploits the possibility to remove the perturbation with the normalization algorithm obtaining, in this way, the desired normal form (5). The interpretation of $\epsilon_{j}$ as a bound for the remainder is clearly related to the well known feature of the quadratic method.

Proof. By induction. If $j=0$, the statement is clearly true by hypothesis, by setting $f^{(0)}:=f$, either in the case $I$ or in the case $I I$. We are supposing here that $\epsilon_{0}$ is small enough in order to satisfy (18). This will be achieved later by a suitable choice of $R_{0}$.
Let us suppose the statement to be valid for $j$. In this way we get

$$
H^{(j+1)} \equiv \exp \left(\mathcal{L}_{\chi^{(j)}}\right) H^{(j)}=h+f^{(j)}+\mathcal{L}_{\chi^{(j)}} h+\sum_{s \geq 1}(s!)^{-1} \mathcal{L}_{\chi^{(j)}}^{s} f^{(j)}+\sum_{s \geq 2}(s!)^{-1} \mathcal{L}_{\chi^{(j)}}^{s} h .
$$

We shall determine $\chi^{(j)}$ in such a way (10) is satisfied so that, by setting

$$
\begin{equation*}
f^{(j+1)}:=\sum_{s \geq 1} \frac{1}{s!} \mathcal{L}_{\chi^{(j)}}^{s} f^{(j)}+\sum_{s \geq 2} \frac{1}{s!} \mathcal{L}_{\chi^{(j)}}^{s} h \stackrel{(10)}{=} \sum_{s \geq 1} \frac{s}{(s+1)!} \mathcal{L}_{\chi^{(j)}}^{s} f^{(j)}, \tag{20}
\end{equation*}
$$

one has $H^{(j+1)}=h+f^{(j+1)}$.
It is immediate from (13) that $\chi^{(j)}$ has the same null Taylor coefficients as $f^{(j)}$. Hence if $f^{(j)}$ is $(Q x L y)$ then $\chi^{(j)}$ is also. It is easy to check by induction that this implies that $\mathcal{L}_{\chi^{(j)}}^{s} f^{(j)}$ is $(Q x L y)$ for all $s$, then $f^{(j+1)}$ is $(Q x L y)$. Similarly, equation (15) implies that if $f^{(j)}$ is $(L y)$ then $\chi^{(j)}$ is also. This implies that $\mathcal{L}_{\chi^{(j)}}^{s} f^{(j)}$ is $(L y)$ for all $s$, hence $f^{(j+1)}$ is $(L y)$. This completes the formal part. In particular, by induction, $f^{(j)}$ is $(L y)$ for all $j$, as claimed in Remark 1.3.
Let us now discuss the quantitative estimate on $f^{(j)}$ in the case $a>0$. By Propositions 2.2, 2.3 and the inductive hypothesis, one gets

$$
\begin{equation*}
\left\|\mathcal{L}_{\chi^{(j)}}^{s} f^{(j)}\right\|_{\left(1-2 d_{j}\right) R_{j}} \leq s!\Theta^{s} \epsilon_{j} e^{-a t}, \quad \Theta:=\frac{e^{2} C_{1}}{a R_{*}^{2} d_{j}^{2 n+4}} \epsilon_{j} \tag{21}
\end{equation*}
$$

Setting $K:=2 n e^{2} C_{1} R_{*}^{-2}$ and $\sigma:=2 n+5$, we have that

$$
\begin{equation*}
2 n \Theta=\left(K \epsilon_{j} a^{-1} d_{j}^{-\sigma}\right) d_{j} \leq d_{j} \tag{22}
\end{equation*}
$$

as $\epsilon_{j+1} / \epsilon_{j}<1$ by Lemma 3.1. Hence, $\Theta<1 / 2$ and the series defined in (20) is convergent, furthermore

$$
\begin{equation*}
e^{a t}\left\|f^{(j+1)}\right\|_{R_{j+1}} \leq \epsilon_{j} \sum_{s \geq 1} \Theta^{s} \leq 2 n \Theta \epsilon_{j} \stackrel{(22)}{\leq} K a^{-1} d_{j}^{-\sigma} \epsilon_{j}^{2} \stackrel{(17)}{=} \epsilon_{j+1} \tag{23}
\end{equation*}
$$

which completes the inductive step. The condition (18) in this case reads as

$$
\begin{equation*}
\epsilon_{0} \leq a R_{0}^{2}(2 \pi)^{-\sigma}\left(8 n e^{2} C_{1}\right)^{-1} \tag{24}
\end{equation*}
$$

On the other hand, from the analyticity of $f$, we get $\left|f_{\alpha, \beta}(t)\right| \leq M_{f} R^{-|\alpha+\beta|} \leq M_{f} R_{0}^{-|\alpha+\beta| / 16}$, as $R_{0} \leq R^{16}$ by hypothesis. By using the first of (9) we get $\|f\|_{R_{0}} \leq M_{f} \sum_{(\alpha, \beta) \in \mathbb{N}^{2 n}} R_{0}^{(15 / 16)|\alpha+\beta|} \leq$ $2 n e^{(2 n-1)} M_{f} R_{0}^{135 / 64}=: \epsilon_{0}$. Replacing the latter in (24), the condition on $R_{0}$ described in the statement of Theorem 1.1 is meant to be completed with the following one

$$
\begin{equation*}
R_{0} \leq\left[a /\left(16(2 \pi)^{\sigma} e^{2 n+1} n^{2} C_{1} M_{f}\right)\right]^{64 / 7} \tag{25}
\end{equation*}
$$

The case $a=0$ is analogous: it is sufficient to replace $C_{1}$ with $C_{2}$, remove the term $e^{ \pm a t}$ from the statement, (21) and (23), then replace $a$ with 1 from (21) to (24), where now $\sigma=n+\tau+5$. The only substantial difference consists in the sum obtained from (9), which is slightly improved, since $f$ linear in $y$. We have in this case $\|f\|_{R_{0}} \leq n^{2} e^{n-1} M_{f} R_{0}^{75 / 32}=: \epsilon_{0}$ leading to

$$
\begin{equation*}
R_{0} \leq\left[8(2 \pi)^{\sigma} e^{n+1} n^{3} C_{2} M_{f}\right]^{-32 / 11} \tag{26}
\end{equation*}
$$

### 3.3 Bounds on the coordinate transformation

Lemma 3.4. The transformation of coordinates defined by the limit (6) satisfies

$$
\begin{equation*}
\left|x^{(\infty)}-x\right|,\left|y^{(\infty)}-y\right|,\left|\eta^{(\infty)}-\eta\right| \leq R_{0} / 6 \tag{27}
\end{equation*}
$$

in particular, it defines an analytic map $\mathcal{M}: \mathcal{D}_{R_{*}} \rightarrow \mathcal{D}_{R_{0}}$ and $H^{(\infty)}:=\mathcal{M} H$ is an analytic function on $\mathcal{D}_{R_{*}}$.

Proof. We will discuss the case $a>0$. The case $a=0$ is straightforward simply replacing $C_{1}$ with $C_{2}, a$ with 1 and changing the value of $\sigma$, where necessary.
Let us start from the variable $x$. Note that, by Proposition 2.3, one has $\left\|\mathcal{L}_{\chi^{(j)}}^{s} x_{l}^{(j+1)}\right\|_{\left(1-2 d_{j}\right) R_{j}} \leq$ $s!\Theta^{s} R_{0}$ for all $l=1, \ldots, n$. Hence we have, by (22)

$$
\left|x^{(j+1)}-x^{(j)}\right| \leq n \max _{l=1, \ldots, n} \sum_{s \geq 1} \frac{1}{s!}\left\|\mathcal{L}_{\chi^{(j)}}^{s} x_{l}^{(j+1)}\right\|_{\left(1-2 d_{j}\right) R_{j}} \leq 2 n R_{0} \Theta \leq R_{0} d_{j}
$$

In this way $\left|x^{(\infty)}-x\right| \leq \sum_{j \geq 0}\left|x^{(j+1)}-x^{(j)}\right|$ converges by (19). The procedure for $y$ is analogous. As for the third of (27), it is necessary to observe that $\mathcal{L}_{\chi^{(j)}} \eta=-\partial_{t} \chi^{(j)}$. Hence, by (16) and the second of $(11)$, one has $\left\|\mathcal{L}_{\chi^{(j)}}^{s} \eta\right\|_{\left(1-2 R_{j}\right)} \leq e^{-2} s!\Theta^{s-1}\left(R_{*}^{2} e^{-2} \Theta\right) \leq s!\Theta^{s} R_{0}$, hence $\left|\eta^{(j+1)}-\eta^{(j)}\right| \leq$ $2 n R_{0} \Theta \leq R_{0} d_{j}$.
The bounds (27) ensure that points in $\mathcal{D}_{R_{*}}$ are mapped within $\mathcal{D}_{R_{0}}$ where $R_{*}=R_{0} / 2$. Furthermore, the absolute convergence of the above described series, ensured by (19), guarantees the uniform convergence in every compact subset of $\mathcal{D}_{R_{*}}$ and the analyticity of $\mathcal{M}$, and then of $H^{(\infty)}$, follows from the theorem of Weierstraß, see e.g. [Det65].

## Appendix. Proof of Proposition 2.1

First of all, recall $\sum_{|\nu| \geq N}|\nu|^{\mu} \mathcal{R}^{|\nu|}=\sum_{l \geq N}\binom{l+m-1}{m-1} l^{\mu} \mathcal{R}^{l}$. Now note that $\log \prod_{j=1}^{m-1}(l+j) \leq$ $\int_{1}^{m} \log (l+x) d x=1-m+\log \left[(m+l)^{(m+l)}(1+l)^{-(1+l)}\right]$ hence $(m-1)!\binom{l+m-1}{m-1}=\prod_{j=1}^{m-1}(l+j) \leq$ $e^{m-1}(m+l)^{(m+l)}(1+l)^{-(1+l)} \leq e^{2 m-2}(m+l)^{(m+\mu)}$. This yields

$$
\begin{equation*}
\sum_{|\nu| \geq N}|\nu|^{\mu} \mathcal{R}^{|\nu|} \leq\left[e^{2 m-2} /(m-1)!\right] \sum_{l \geq N}(m+l)^{(m+\mu)} \mathcal{R}^{l} \tag{28}
\end{equation*}
$$

On the other hand, the function $h(x):=(m+x)^{\kappa} \mathcal{R}^{x / 4}$ has a maximum in $x=0$ (in the nonnegative semi-axis) if $\mathcal{R} \leq \exp (-4 \kappa / m)$ and in $x^{*}:=-m-4 \kappa / \log \mathcal{R}$ otherwise. Hence, from (28) with $\mu=0$ we have $\sum_{|\nu| \geq N} \mathcal{R}^{|\nu|} \leq[(m-1)!]^{-1} m^{m} e^{2 m-2} \sum_{l \geq N} \mathcal{R}^{(3 / 4) l}$ which gives the first of (9) by using the inequality $m^{m} \leq e^{m-1} m$ ! and recalling $\mathcal{R} \leq e^{-4}$.
Now set $\mathcal{R}=1-\delta$. By hypothesis $\mathcal{R}>e^{-4}$, hence $(m+l)^{(m+\mu)}(1-\delta)^{l / 4} \leq(1-\delta)^{-m / 2}(-2(m+$ $\mu) / \log (1-\delta))^{(m+\mu)}$. By substituting the latter in (28) with $N=0$, then using the inequalities $-\log (1-\delta) \geq \delta$ and $\left[1-(1-\delta)^{3 / 4}\right] \geq \delta / 2$ as $\delta \leq 1 / 2$, the second of (9) easily follows.

## Acknowledgements

The first author is grateful to Prof. Dario Bambusi for remarkable discussions on this problem.

## References

[Ber09] V. Berdichevsky. Variational Principles of Continuum Mechanics: I. Fundamentals. Interaction of Mechanics and Mathematics. Springer Berlin Heidelberg, 2009.
[Bou13] A. Bounemoura. Effective stability for slow time-dependent near-integrable hamiltonians and application. C. R. Math. Acad. Sci. Paris, 351(17-18):673-676, 2013.
[CG94] L. Chierchia and G. Gallavotti. Drift and diffusion in phase space. Ann. Inst. H. Poincaré Phys. Théor., 60(1):144, 1994.
[Det65] J.W. Dettman. Applied Complex Variables. A Series of advanced mathematics texts. Dover Publications, 1965.
[FW14] A. Fortunati and S. Wiggins. Normal form and Nekhoroshev stability for nearly integrable Hamiltonian systems with unconditionally slow aperiodic time dependence. Regul. Chaotic Dyn., 19(3):363-373, 2014.
[FW15a] A. Fortunati and S. Wiggins. A Kolmogorov theorem for nearly-integrable Poisson systems with asymptotically decaying time-dependent perturbation. Regul. Chaotic Dyn., 20(4):476-485, 2015.
[FW15b] A. Fortunati and S. Wiggins. Normal forms à la Moser for aperiodically time-dependent Hamiltonians in the vicinity of a hyperbolic equilibrium. Accepted for the publication on Discr. Cont. Dyn. Sys. -S., 2015.
[FW15c] A. Fortunati and S. Wiggins. Negligibility of small divisor effects in the normal form theory for nearly-integrable hamiltonians with decaying non-autonomous perturbations. http://arxiv.org/abs/1509.02119, 2015.
[Gio] A. Giorgilli. On a Theorem of Lyapounov. Rendiconti dell'Istituto Lombardo Accademia di Scienze e Lettere, Classe di Scienze Matematiche e Naturali, In Press.
[Gio03] A. Giorgilli. Exponential stability of Hamiltonian systems. In Dynamical systems. Part I, Pubbl. Cent. Ric. Mat. Ennio Giorgi, pages 87-198. Scuola Norm. Sup., Pisa, 2003.
[GZ92] A. Giorgilli and E. Zehnder. Exponential stability for time dependent potentials. Z. Angew. Math. Phys., 43(5):827-855, 1992.
[Lya92] A.M. Lyapunov. General Problem of the Stability Of Motion. Control Theory and Applications Series. Taylor \& Francis, 1992.
[Mos56] J. Moser. The analytic invariants of an area-preserving mapping near a hyperbolic fixed point. Comm. Pure Appl. Math., 9:673-692, 1956.
[Poi79] H. Poincaré. Sur les propriétés des fonctions définies par les équations aux différences partielles. Paris. Université. Faculté des sciences. Theses. Gauthier-Villars, 1879.
[Pus74] L. D. Pustyl'nikov. Stable and oscillating motions in nonautonomous dynamical systems. A generalization of C. L. Siegel's theorem to the nonautonomous case. Mat. $S b$. (N.S.), 94(136):407-429, 495, 1974.
[Sie42] C. L. Siegel. Iteration of analytic functions. Ann. of Math. (2), 43:607-612, 1942.
[Sie52] C. L. Siegel. Über die Normalform analytischer Differentialgleichungen in der Nähe einer Gleichgewichtslösung. Nachr. Akad. Wiss. Göttingen. Math.-Phys. Kl. Math.-Phys.-Chem. Abt., 1952:21-30, 1952.
[Sie54] C. L. Siegel. Über die Existenz einer Normalform analytischer Hamiltonscher Differentialgleichungen in der Nähe einer Gleichgewichtslösung. Math. Ann., 128:144-170, 1954.


[^0]:    ${ }^{*}$ This research was supported by ONR Grant No. N00014-01-1-0769 and MINECO: ICMAT Severo Ochoa project SEV-2011-0087.
    ${ }^{(a)}$ E-mail: alessandro.fortunati@bristol.ac.uk
    ${ }^{(b)}$ E-mail: s.wiggins@bristol.ac.uk

[^1]:    ${ }^{1}$ I.e. such that the corresponding canonical equations are integrable.
    ${ }^{2}$ With periodic time dependence.

[^2]:    ${ }^{3}$ The exponential decay, see (3), is chosen for simplicity of discussion. The only necessary assumption is the summability in $t$ of the perturbing function over the non-negative real semi-axis, see [FW15c].
    ${ }^{4}$ It is understood that $x^{\alpha} y^{\beta}:=x_{1}^{\alpha_{1}} \cdot \ldots \cdot x_{n}^{\alpha_{n}} \cdot y_{1}^{\beta_{1}} \cdot \ldots \cdot y_{n}^{\beta_{n}}$.
    ${ }^{5}$ The interval $a \in[0,1)$ is a compact way to denote either the time decay $a \in(0,1)$ or the boundedness $a=0$. As in our previous paper we recall that we are interested in the case of small $a$ (slow decay) and the upper bound $a=1$ is set for simplicity. On the other hand, it is easy to realise that the case $a \geq 1$ is straightforward.

