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ON OSCILLATIONS IN A SYSTEM WITH A PIECEWISE SMOOTH COEFFICIENT

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Abstract
In this paper a straightforward perturbation procedure will be presented for a class of ordinary differential equations with piecewise smooth terms in the equations. From the constructed approximations of the solutions a map can be derived from which the existence and the stability of time-periodic solutions can be determined.

1 Introduction

Ordinary differential equations as model equation for piecewise smooth systems occur in a wide variety of applications. For dry friction problems the equation

\[ m\ddot{x} + c \dot{x} = F_0 \text{sgn}(v_0 - \dot{x}) \]  

(1)

is frequently used. In Eq. (1) \( m, c, F_0, \) and \( v_0 \) are constants and \( \text{sgn} \) is the signum function defined by \( \text{sgn}(u) = 1 \) for \( u > 0 \), \( 0 \) for \( u = 0 \), and \( -1 \) for \( u < 0 \). For oscillation problems with a symmetrical restoring force of constant magnitude the equation

\[ \ddot{x} + \text{sgn}(x) = 0 \]  

(2)

can be used. The Eqs. (1) and (2), and many similar equations can for instance be found in Jordan and Smith [1]. Another example of an equation
for a piecewise smooth system is

$$\ddot{x} + c\dot{x} + \mu x^+ - \nu x^- = f(t, x, \dot{x}),$$  \hspace{1cm} (3)$$

which has been introduced by Lazer and McKenna [2] as a simple model equation for the vertical oscillations of a long-span suspension bridge. In Eq. (3) $c$, $\mu$, and $\nu$ are positive constants, $f(t, x, \dot{x})$ is an external force function, and $x^+ = \max\{x, 0\}$ is the positive part of $x$, and $x^- = \max\{-x, 0\}$ its negative part. Recently Hogan [3] studied the equation

$$\ddot{x} + \varepsilon(|x| - 1)\dot{x} + x = 0,$$  \hspace{1cm} (4)$$

where $\varepsilon$ is a small positive constant, that is, $0 < \varepsilon < 1$. Eq. (4) is related to the equation

$$\ddot{x} + \varepsilon(|\dot{x}| - 1)\dot{x} + x = 0,$$  \hspace{1cm} (5)$$

which has been studied in Jordan and Smith [1]. In fact by differentiating Eq. (5) with respect to $t$ (and by rescaling $x$ by a factor 2) an equation like Eq. (4) is obtained. In [1] the averaging technique and the harmonic balance method have been used to construct approximations of the periodic solution of Eq. (5) and in [3] the two-time-scales perturbation method has been applied to Eq. (4). To prove the existence of a unique, nontrivial periodic solution a Liénard theorem has been used in [3]. Usually it is difficult to find (if available) the appropriate Liénard theorem for a given differential equation. In this paper it will be shown that for perturbed equations (like (1)-(5)) a straightforward perturbation procedure can be used to construct approximations of the solutions. From these approximations a map can be defined. By using this map the existence and the stability of time-period solutions can be established. As a prototype of problem Eq. (4) will be treated in detail.
2 Construction of an approximation of the solution

In this section an approximation of the solution of Eq. (4) will be constructed. Since Eq. (4) is autonomous it can be assumed without loss of generality that \( x(0) = 0 \) and \( \dot{x}(0) = A > 0 \). So, the following problem has to be considered (as long as \( x(t) \geq 0 \) for \( t > 0 \)):

\[
\begin{align*}
\ddot{x} + \varepsilon (x - 1) \dot{x} + x &= 0, \quad t > 0, \\
x(0) &= 0, \quad \dot{x}(0) = A > 0,
\end{align*}
\] (6)

where \( 0 < \varepsilon \ll 1 \), and where \( A \) is an \( \varepsilon \)-independent constant. To solve the initial-value problem for (6) approximately the following expansion in \( \varepsilon \) for \( x(t) \), that is,

\[
x_0(t) + \varepsilon x_1(t) + \varepsilon^2 x_2(t) + \ldots
\] (7)

is substituted into (6). By taking together terms of equal powers in \( \varepsilon \) the following problems are then obtained:

\[
O(1)\text{-problem : } \ddot{x}_0 + x_0 = 0, \quad t > 0, \\
x_0(0) = 0, \quad \dot{x}_0(0) = A,
\] (8)

\[
O(\varepsilon)\text{-problem : } \ddot{x}_1 + x_1 = -(x_0 - 1) \dot{x}_0, \quad t > 0, \\
x_1(0) = \dot{x}_1(0) = 0,
\] (9)

and so on. The \( O(1) \)-problem (8) can readily be solved, yielding

\[
x_0(t) = A \sin(t).
\] (10)

Substituting (10) into (9) the \( O(\varepsilon) \)-problem (9) then also can be solved, yielding

\[
x_1(t) = -\frac{A^2}{3} \sin(t) + \frac{A^2}{6} \sin(2t) + \frac{A}{2} t \sin(t).
\] (11)

Now it should be observed that the approximation \( \bar{x}(t) = x_0(t) + \varepsilon x_1(t) \) satisfies the original differential equation (6) and the initial values up to order \( \varepsilon^2 \), that is,
\[ \ddot{x} + \varepsilon (\ddot{x} - 1) \dot{x} + \dddot{x} = \]
\[ \dddot{x}_0 + x_0 + \varepsilon (\dddot{x}_1 + x_1 + (x_0 - 1) \dddot{x}_0) + \varepsilon^2 ((x_0 - 1) \dddot{x}_1 + \dddot{x}_0 \dot{x}_1) + \varepsilon^3 x_1 \dot{x}_1 = O(\varepsilon^2 t) + O(\varepsilon^3 t^2), \]
and \( \dddot{x}(0) = 0, \quad \dot{x}(0) = A \). \quad (12)

By writing down the equivalent integral equations for the initial value problems (6) and (12), by subtracting the so-obtained integral equations, by using the Lipschitz continuity of the nonlinearities, and by applying Gronwall’s inequality it can be shown elementarily (see for instance chapter 1 of Verhulst [4]) that (as long as \( x(t) \geq 0 \))

\[ |x(t) - \dddot{x}(t)| = O(\varepsilon^2 t^2). \quad (13) \]

The estimate (13) implies that (as long as \( x(t) \geq 0 \)) \( x(t) = \dddot{x}(t) + O(\varepsilon^2) = x_0(t) + \varepsilon x_1(t) + O(\varepsilon^2) \) for times \( t \) of order 1, where \( x_0(t) \) and \( x_1(t) \) are given by (10) and (11) respectively. Now let \( \tilde{t} > 0 \) be the shortest time for which \( x(t) = x_0(t) + \varepsilon x_1(t) + O(\varepsilon^2) \) becomes zero, and assume that \( \tilde{t} \) can be approximated by \( \tilde{t}_0 + \varepsilon \tilde{t}_1 + \varepsilon^2 \tilde{t}_2 + \ldots \). By substituting this approximation for \( \tilde{t} \) into \( x(\tilde{t}) = x_0(\tilde{t}) + \varepsilon x_1(\tilde{t}) + O(\varepsilon^2) = 0 \) it easily follows that \( \tilde{t}_0 = \pi, \) and \( \tilde{t}_1 = 0, \) and so \( \tilde{t} = \pi + O(\varepsilon^2). \) It then follows from \( x(t) = x_0(t) + \varepsilon x_1(t) + O(\varepsilon^2) \) that (for \( \varepsilon \) sufficiently small) \( \dot{x}(\tilde{t}) < 0. \) So, for \( t > \tilde{t} \) the following initial value problem has to solved (until \( x \) becomes zero again):

\[ \ddot{x} + \varepsilon (-x - 1) \dot{x} + x = 0, \quad t > \tilde{t}, \]
\[ x(\tilde{t}) = 0, \quad \dot{x}(\tilde{t}) = -A + \varepsilon A \left( \frac{2}{3} A - \frac{\pi}{2} \right) + O(\varepsilon^2), \quad (14) \]

where \( \tilde{t} = \pi + O(\varepsilon^2). \) To solve this initial value problem (14) approximately the function \( x(t) \) is again approximated by the expansion (7). This expansion is submitted into (14), and terms of equal powers in \( \varepsilon \) are taken together,
yielding as

\[ O(1)\)-problem : \quad \ddot{x}_0 + x_0 = 0 , \quad t > \tilde{t} , \]
\[ x_0(\tilde{t}) = 0 , \quad \dot{x}_0(\tilde{t}) = -A , \quad (15) \]

\[ O(\varepsilon)\)-problem : \quad \ddot{x}_1 + x_1 = (x_0 + 1)\dot{x}_0 , \quad t > \tilde{t} , \]
\[ x_1(\tilde{t}) = 0 , \quad \dot{x}_1(\tilde{t}) = A\left(\frac{2}{3}A - \frac{\pi}{2}\right) , \quad (16) \]

and so on. Again the initial value problems (15), (16), and so on can be solved consecutively, yielding

\[ x_0(t) = -A \sin(t - \tilde{t}) , \quad (17) \]

\[ x_1(t) = A \left( A - \frac{\pi}{2} \right) \sin(t - \tilde{t}) - \frac{A^2}{6} \sin(2(t - \tilde{t})) - \frac{A}{2}(t - \tilde{t}) \sin(t - \tilde{t}) , \quad (18) \]

and so on. Again it can be shown elementarily that the approximation \( \bar{x}(t) = x_0(t) + \varepsilon x_1(t) \) (where \( x_0 \) and \( x_1 \) are given by (17) and (18) respectively) satisfies the following estimate: \( |x(t) - \bar{x}(t)| = O(\varepsilon^2) \) for times \( t \) of order 1 (as long as \( x(t) \leq 0 \)). Now let \( \hat{t} > \tilde{t} \) be the shortest time for which \( x(t) = x_0(t) + \varepsilon x_1(t) + O(\varepsilon^2) \) becomes zero again, and assume that \( \hat{t} \) can be approximated by \( \hat{t}_0 + \varepsilon \hat{t}_1 + \varepsilon^2 \hat{t}_2 + .... \). By substituting this approximation for \( \hat{t} \) into \( x(\hat{t}) = x_0(\hat{t}) + \varepsilon x_1(\hat{t}) + O(\varepsilon^2) = 0 \), where \( x_0 \) and \( x_1 \) are given by (17) and (18) respectively, it easily follows that \( \hat{t}_0 = 2\pi \), and \( \hat{t}_1 = 0 \), and so \( \hat{t} = 2\pi + O(\varepsilon^2) \). Now \( \dot{x}(\hat{t}) \) can also be computed, yielding

\[ \dot{x}(\hat{t}) = A + \varepsilon A(-\frac{4}{3}A + \pi) + O(\varepsilon^2) . \quad (19) \]

For \( \varepsilon \) sufficiently small it follows from (19) that \( \dot{x}(\hat{t}) > 0 \). So, for \( t > \hat{t} \) Eq. (6) has to be solved again subject to \( x(\hat{t}) = 0 \) and \( \dot{x}(\hat{t}) \) as given by (19). Since \( A > 0 \) is arbitrary it can be concluded that in fact all initial value problems for Eq. (4) have been solved already (up to \( O(\varepsilon^2) \)) for times of order 1. Comparing \( \dot{x}(0) = A \) with \( \dot{x}(\hat{t}) \) as given by (19) it can also be concluded that for \( A > \frac{3}{4}\pi \) (and \( \varepsilon \) sufficiently small and positive) the amplitude of the oscillation will decrease, and that for \( A < \frac{3}{4}\pi \) the amplitude will increase.
This suggests the existence of (at least) one nontrivial periodic solution (or closed orbit in the \((x, \dot{x})\) phase plane). In the next section it will be proven that there exists a unique, nontrivial periodic solution for Eq. (4).

### 3 The existence of time-periodic solutions

To prove the existence of a time-periodic solution of Eq. (4) a map \(Q\) (actually the Poincaré-return map on the section \(x = 0\), and \(\dot{x} > 0\) in the \((x, \dot{x})\) phase-plane) will be considered. Based upon the error estimate (13) and upon (19) this map \(Q\) is defined by (for \(n = 1, 2, 3, \ldots\)):

\[
A_n = A_{n-1} + \varepsilon A_{n-1} \left( -\frac{4}{3} A_{n-1} + \pi \right) + O(\varepsilon^2 n),
\]

where \(A_n\) is the value of \(\dot{x}(t)\) with \(t\) such that the solution \(x(t)\) returned to the Poincaré section (that is, \(x = 0\) and \(\dot{x} > 0\)) for the \(n\)th time. A map \(Q: A \to Q(A) \iff A_n = Q(A_{n-1})\) has now be defined. By neglecting terms of \(O(\varepsilon^2 n)\) in (20) a new map \(P\) is defined: \(\tilde{A} \to P(\tilde{A}) \iff \tilde{A}_n = P(\tilde{A}_{n-1})\). The maps \(P\) and \(Q\) will now be studied, and it will be shown that:

(i) if \(|A_0 - \tilde{A}_0| = O(\varepsilon)\) for \(\varepsilon \downarrow 0\), then \(|A_n - \tilde{A}_n| = O(\varepsilon)\) for \(n = O(\frac{1}{\sqrt{\varepsilon}})\),

that is, for \(n \sim \frac{1}{\sqrt{\varepsilon}}\) and \(\varepsilon \downarrow 0\), \(A_n\) and \(\tilde{A}_n\) remain “\(\varepsilon\)-close”;  

(ii) the map \(P\) has a unique, hyperbolic fixed point \(\tilde{A} = \frac{3\pi}{4}\), which is asymptotically stable;

(iii) there exists an \(\varepsilon_0 > 0\) such that for all \(0 < \varepsilon \leq \varepsilon_0\) the map \(Q\) has a unique, hyperbolic fixed point \(A = \frac{3\pi}{4} + O(\varepsilon)\) with the same stability property as the fixed point \(A = \frac{3\pi}{4}\) of the map \(P\).

**Proof of (i).** From \(|A_0 - \tilde{A}_0| = O(\varepsilon)\) for \(\varepsilon \downarrow 0\) it follows that there exists a positive constant \(M_0\) such that \(|A_0 - \tilde{A}_0| = M_0\varepsilon\). Further, \(|A_n - \tilde{A}_n| = |P(A_{n-1}) - P(\tilde{A}_{n-1}) + O(\varepsilon^2 n)| \leq |P(A_{n-1}) - P(\tilde{A}_{n-1})| + M_1\varepsilon^2 n\), where \(M_1\) is a positive constant. We also have \(|P(A_{n-1}) - P(\tilde{A}_{n-1})| \leq L|A_{n-1} - \tilde{A}_{n-1}|\),
where $L$ is a Lipschitz-constant with $L = 1 + \varepsilon M_2$ in which $M_2$ is a positive constant. And so, $|A_n - \tilde{A}_n| \leq (1 + \varepsilon M_2)|A_{n-1} - \tilde{A}_{n-1}| + M_1 \varepsilon^2 n \leq \ldots \leq \varepsilon (M_0 + \varepsilon n^2 M_1) e^{\varepsilon n M_2}$. Thus for $n = O(\frac{1}{\sqrt{\varepsilon}})$ it follows that $|A_n - \tilde{A}_n| = O(\varepsilon)$.

Proof of (ii). The fixed points of the map $P$ follow from $\tilde{A}_n = P(\tilde{A}_{n-1})$ for $n \to \infty$, or equivalently from $\tilde{A} = \tilde{A} + \varepsilon \tilde{A}(-\frac{1}{3} \tilde{A} + \pi) \Leftrightarrow \tilde{A}(-\frac{1}{3} \tilde{A} + \pi) = 0$. Therefore for $\tilde{A} > 0$ the unique fixed point is $\tilde{A} = \frac{4}{3} \pi$. If the linearized map of $P$ around this fixed point has no eigenvalues of unit modulus, then this fixed point is hyperbolic. Let $DP$ be this linearized map; then $DP = 1 - \varepsilon \pi$. Obviously the eigenvalue $\lambda$ of $DP$ is $1 - \varepsilon \pi$. Since $0 < \varepsilon \ll 1$ it follows that $|\lambda| < 1$, and so the fixed point is hyperbolic and stable.

Proof of (iii). It follows from (ii) that $I - DP$ is invertible. Therefore from (i), the implicit function theorem, and the linearization theorem of Hartman-Grobman, it follows (see also chapter 1 in Guckenheimer and Holmes [5] or Van Horssen [6]), that the map $Q$ continues to have a unique fixed point, which is $\varepsilon$-close to $\frac{3\pi}{4}$. Furthermore, the eigenvalue of $DQ$ ($DQ$ is the linearized map of $Q$ around its fixed point) depends continuously on $\varepsilon$, and is (up to $O(\varepsilon^{3/2})$) equal to the eigenvalue of $DP$. Therefore for sufficiently small $\varepsilon$ the two fixed points have the same stability properties.

So far it can be concluded that a unique, asymptotically stable, nontrivial periodic solution exists for Eq. (4). The period $T$ of this unique, nontrivial periodic solution can be approximated by $2\pi + O(\varepsilon^2)$, and the periodic solution itself can be approximated by

$$\frac{3\pi}{4} \sin(t) + \varepsilon \left( -\frac{3\pi^2}{16} \sin(t) + \frac{3\pi^2}{32} \sin(2t) + \frac{3\pi}{8} \sin(t) \right) + O(\varepsilon^2)$$

for $0 \leq t \leq \tilde{t}$, and by

$$\frac{3\pi}{4} \sin(t) + \varepsilon \left( -\frac{9\pi^2}{16} \sin(t) - \frac{3\pi^2}{32} \sin(2t) + \frac{3\pi}{8} \sin(t) \right) + O(\varepsilon^2)$$

for $\tilde{t} \leq t \leq T$, where $\tilde{t}$ is the time as defined before and $\tilde{t}$ can be approximated by $\pi + O(\varepsilon^2)$. Higher order approximations of the period $T$ and of the periodic...
solution can be obtained by solving the $O(\varepsilon^n)$-problems with $n \geq 2$ as defined in section 2 of this paper.

4 Conclusions and remarks

In this paper it has been shown how in a straightforward way approximations of the solutions can be obtained for a class of ordinary differential equations with piecewise smooth terms in the equations. By using these approximations a map can be constructed from which the existence and the stability of time-periodic solutions can be determined. As example the equation

$$\ddot{x} + \varepsilon(|x| - 1)\dot{x} + x = 0$$

has been studied in detail. Not only approximations of the solutions of this equation have been constructed but also order $\varepsilon^2$ accurate approximations of the unique and stable periodic solution have been presented. The perturbation procedure as presented in this paper can be applied to a large class of ordinary differential equations with piecewise smooth terms in the equations. In a forthcoming paper the equation

$$\ddot{x} + \mu x^+ - \nu x^- = \varepsilon(\dot{x} - \dot{x}^3) \tag{21}$$

will be studied. In this equation $\mu, \nu,$ and $\varepsilon$ are positive constants with $0 < \varepsilon \ll 1$, and $x^+ = \max\{x, 0\}$, and $x^- = \max\{-x, 0\}$. This equation serves as a simple model equation to describe the wind-induced vertical vibrations of a long-span suspension bridge. When $\mu = \nu$ the well-known Rayleigh equation is obtained, which is related to the Van der Pol equation. The presented perturbation procedure can also be applied successfully to Eq. (21), and it can be shown that also Eq. (21) has a unique and stable periodic solution.
References


