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# ON THE EXISTENCE OF OSCILLATORY SOLUTIONS IN THE WEISBUCH-SALOMON-ATLAN MODEL FOR THE BELOUSOV-ZHABOTINSKIJ REACTION 

Valter Šeda

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The Belousov-Zhabotinskij reaction is an oscillating oxidation reaction. For the kinetics of this reaction two models have been established. The first of them has been investigated in [4] where the existence of a periodic solution in the model has been proved. As was pointed out, all solutions in this model are oscillatory. In [7] G. Weisbuch, J. Salomon and H. Atlan have proposed a new model and have examined the stability of this model. They have shown numerically that this system shows oscillations, too. In this paper the stability properties of the model are reexamined and completed, the existence and properties of oscillatory solutions are investigated and a periodic solution is found.

1. The model in question is

$$
\begin{align*}
\dot{X} & =-5 K_{1} X+K_{5} C  \tag{1}\\
\dot{U} & =K_{1} X+K_{3} U-K_{2} X U-K_{4} U^{2} \\
\dot{C} & =2 K_{3} U-4 K_{5} C
\end{align*}
$$

where $K_{1}-K_{5}$ are positive real parameters representing kinetic constants and $X$, $U, C$ are concentrations, and hence, nonnegative.

The system (1) has two equilibrium points

$$
\begin{equation*}
a_{0}=(0,0,0), \quad a_{1}=\left(X_{0}, U_{0}, C_{0}\right) \tag{2}
\end{equation*}
$$

where

$$
\begin{align*}
X_{0} & =\frac{11 K_{3}^{2}}{10\left(K_{2} K_{3}+10 K_{1} K_{4}\right)}, \quad U_{0}=\frac{11 K_{1} K_{3}}{K_{2} K_{3}+10 K_{1} K_{4}},  \tag{3}\\
C_{0} & =\frac{11 K_{1} K_{3}^{2}}{2 K_{5}\left(K_{2} K_{3}+10 K_{1} K_{4}\right)}
\end{align*}
$$

and thus

$$
\begin{equation*}
X_{0}=\frac{K_{3}}{10 K_{1}} U_{0}=\frac{K_{5}}{5 K_{1}} C_{0}, \quad C_{0}=\frac{K_{3}}{2 K_{5}} U_{0} \tag{4}
\end{equation*}
$$

First a general property of the system (1) will be derived. Denote $\mathscr{P}=\{(X, U, C)$ : $: X \geqq 0, U \geqq 0, C \geqq 0\}$ and let $\mathscr{P}^{0}$ be the interior of $\mathscr{P}$. Since $a_{0}$ is a non-egress point and all the other points of quarterplanes $X=0(U \geqq 0, C \geqq 0), U=0(X \geqq 0$, $C \geqq 0)$, and $C=0(X \geqq 0, U \geqq 0)$ are strict ingrees points of $P^{0}$ with respect to (1), we obtain on the basis of Lemma 8.1, [3], p. 53,

Theorem 1. The system (1) satisfies Hypothesis $V$ from paper [2], p. 31, i.e., for each solution $(X, U, C)$ of $(1)$ it holds that if there is a $t_{0}$ such that $\left(X\left(t_{0}\right), U\left(t_{0}\right)\right.$, $\left.C\left(t_{0}\right)\right) \in \mathscr{P}$, then $(X(t), U(t), C(t)) \in \mathscr{P}$ for all $t \geqq t_{0}$ (from the interval of existence).

Now the stability of critical points $a_{0}, a_{1}$ will be investigated. To this aim let us introduce new variables $x, u, c$ by

$$
\begin{align*}
& X=X_{0}+x,  \tag{5}\\
& U=U_{0}+u, \\
& C=C_{0}+c
\end{align*}
$$

where $\left(X_{0}, U_{0}, C_{0}\right)$ can also stand for $a_{0}$. Then (1) will assume the form

$$
\begin{align*}
& \dot{x}=-5 K_{1} x+K_{5} c,  \tag{6}\\
& \dot{u}=\left(K_{1}-K_{2} U_{0}\right) x+\left(K_{3}-K_{2} X_{0}-2 K_{4} U_{0}\right) u-K_{2} x u-K_{4} u^{2}, \\
& \dot{c}=2 K_{3} u-4 K_{5} c .
\end{align*}
$$

The characteristic equation of the matrix of the system of the first approximation is of the form

$$
\begin{equation*}
\lambda^{3}+a \lambda^{2}+b \lambda+d=0 \tag{7}
\end{equation*}
$$

where

$$
\begin{align*}
a= & 5 K_{1}-K_{3}+4 K_{5}+K_{2} X_{0}+2 K_{4} U_{0}  \tag{8}\\
b= & 20 K_{1} K_{5}-5 K_{1} K_{3}-4 K_{3} K_{5}+\left(5 K_{1} K_{2}+4 K_{2} K_{5}\right) X_{0}+ \\
& +\left(10 K_{1} K_{4}+8 K_{4} K_{5}\right) U_{0} \\
d= & -22 K_{1} K_{3} K_{5}+20 K_{1} K_{2} K_{5} X_{0}+\left(40 K_{1} K_{4} K_{5}+2 K_{2} K_{3} K_{5}\right) U_{0} .
\end{align*}
$$

Thus in the case of the critical point $a_{0}$

$$
\begin{align*}
& a=5 K_{1}-K_{3}+4 K_{5},  \tag{0}\\
& b=20 K_{1} K_{5}-5 K_{1} K_{3}-4 K_{3} K_{5}, \\
& d=-22 K_{1} K_{3} K_{5},
\end{align*}
$$

while for $a_{1}$ we get

$$
\begin{align*}
a= & 5 K_{1}+K_{3} / 10+4 K_{5}+11 K_{1} K_{3} K_{4} /\left(K_{2} K_{3}+10 K_{1} K_{4}\right),  \tag{1}\\
b= & 20 K_{1} K_{5}+(1 / 2) K_{1} K_{3}+(2 / 5) K_{3} K_{5}+K_{1} K_{3} K_{4}\left(55 K_{1}+\right. \\
& \left.+44 K_{5}\right) /\left(K_{2} K_{3}+10 K_{1} K_{4}\right), \\
d= & 22 K_{1} K_{3} K_{5} .
\end{align*}
$$

The next lemma brings a statement on zeros of the polynomial (7) which puts together a new result obtained with a little algebra, a stability criterion from [6], p. 58, and a result on instability from [7], p. 74.

Lemma 1. I. Let $a>0, b>0, d>0$. Then there is a real negative zero of (7) and, if $d-a b<0$, all zeros of this polynomial have negative real part. If $d-$ $-a b \geqq 0$, then there is a pair $\beta \pm i \gamma$ of complex conjugate roots of (7) such that $\operatorname{sgn} \beta=\operatorname{sgn}(d-a b)$.
II. If $d<0$, then there exists a positive root of (7) and if this polynomial possesses a pair $\beta \pm i \gamma$ of complex conjugate roots, then $\operatorname{sgn} \beta=\operatorname{sgn}(d-a b)$.

Let us consider the problem of stability of $a_{0}$. With respect to $\left(8_{0}\right)$, Lemma 1 implies that (7) has a positive root. Further, (1) satisfies the condition (1.8) in [1], p. 316. Hence Theorem 1.2 [1], p. 317 yields

Theorem 2. The equilibrium point $a_{0}$ of the system (1) is unstable.
As to the stability of the critical point $a_{1}$ of the system (1), we see from (5) that it is equivalent to the stability of the equilibrium point $a_{0}$ of (6). Here, by ( $8_{1}$ ), we have to deal with the case $a>0, b>0, d>0$ of coefficients of (7).

Denote

$$
\begin{equation*}
K_{1}\left|K_{2}=k, \quad K_{3}\right| K_{4}=l . \tag{9}
\end{equation*}
$$

When $k=l$, i.e. $K_{1}=k K_{2}, K_{3}=k K_{4}$, we get from ( 81 ) that

$$
\begin{aligned}
& a=k\left(5 K_{2}+1,1 K_{4}\right)+4 K_{5}, \\
& b=k^{2} 6 K_{2} K_{4}+k\left(20 K_{2} K_{5}+4,8 K_{4} K_{5}\right), \\
& d=k^{2} 22 K_{2} K_{4} K_{5} .
\end{aligned}
$$

Then

$$
d-a b=k^{2}\left(22 K_{2} K_{4} K_{5}-24 K_{2} K_{4} K_{5}\right)+\ldots
$$

where all omitted terms are negative. Therefore

$$
d-a b<-2 K_{2} K_{4} K_{5} k^{2}<0
$$

and from Lemma 1 by the Ljapunov Theorem, [6], p. 213 we obtain

Theorem 3. If $K_{1} / K_{2}=K_{3} / K_{4}$, then the equilibrium point $a_{1}$ of the system (1) is exponentially asymptotically stable.

Corollary. In this case the system (1) is a generalized Volterra equation.
Proof. By Theorem 1, the system (1) satisfies Hypothesis V and has at least one asymptotically stable solution which passes through $\mathscr{P}$ and is bounded. By Definition 4, [2], p. 33 this means that (1) is a generalized Volterra equation.
The case $K_{1} / K_{2} \neq K_{3} / K_{4}$ is more interesting. From ( $\left(8_{1}\right)$ it follows that

$$
\begin{equation*}
d-a b=H_{1}-H_{2} \tag{10}
\end{equation*}
$$

where

$$
\begin{align*}
H_{1}= & 16 K_{1} K_{3} K_{5}-2,5 K_{1}^{2} K_{3}-100 K_{1}^{2} K_{5}-0,05 K_{1} K_{3}^{2}-  \tag{11}\\
& -0,04 K_{3}^{2} K_{5}-80 K_{1} K_{5}^{2}-1,6 K_{3} K_{5}^{2}
\end{align*}
$$

and

$$
\begin{align*}
H_{2}= & {\left[1 /\left(K_{2} K_{3}+10 K_{1} K_{4}\right)\right] \cdot\left[660 K_{1}^{2} K_{3} K_{4} K_{5}+176 K_{1} K_{3} K_{4} K_{5}^{2}+\right.}  \tag{12}\\
& +8,8 K_{1} K_{3}^{2} K_{4} K_{5}+11 K_{1}^{2} K_{3}^{2} K_{4}+275 K_{1}^{3} K_{3} K_{4}+ \\
& +605 K_{1}^{3} K_{3}^{2} K_{4}^{2} /\left(K_{2} K_{3}+10 K_{1} K_{4}\right)+ \\
& \left.+484 K_{1}^{2} K_{3}^{2} K_{4}^{2} K_{5} /\left(K_{2} K_{3}+10 K_{1} K_{4}\right)\right] .
\end{align*}
$$

Note that $H_{1}$ depends only on $K_{1}, K_{3}, K_{5}$ and that for fixed $K_{1}, K_{3}, K_{4}, K_{5}$ the expression $H_{1}-H_{2}$ is an increasing function of $K_{2}$. Further, the relations $\lim _{K_{2} \rightarrow 0+}\left(H_{1}-H_{2}\right)<0, \lim _{K_{2} \rightarrow \infty}\left(H_{1}-H_{2}\right)=H_{1}$ are true. Thus, for such $K_{1}, K_{3}, K_{5}$ that $H_{1}>0$ and for an arbitrary $K_{4}$ there exists a unique $K_{2}=K_{2}^{0}$ for which $H_{1}$--$-H_{2}=0$. Similarly $\lim _{K_{4} \rightarrow 0+}\left(H_{1}-H_{2}\right)=H_{1}$ and $\lim _{K_{4} \rightarrow \infty}\left(H_{1}-H_{2}\right)=-\infty$ imply that if $H_{1}>0$ and $K_{2}$ is arbitrary, there is at least one $K_{4}=K_{4}^{0}$ such that $H_{1}-H_{2}=$ $=0$.
A sufficient condition for $H_{1}$ to be positive will be now derived. Similarly as in [7], p. 75 we put

$$
\begin{align*}
& K_{3}=\nu K_{1},  \tag{13}\\
& K_{1}=\mu K_{5} .
\end{align*}
$$

Then

$$
\begin{equation*}
H_{1}=\mu K_{5}^{3}\left(-p \mu^{2}+q \mu-r\right), \tag{14}
\end{equation*}
$$

where

$$
\begin{align*}
& p=2,5 v+0,05 v^{2},  \tag{15}\\
& q=16 v-100-0,04 v^{2}, \\
& r=1,6 v+80 .
\end{align*}
$$

Let $v$ be such that

$$
\begin{equation*}
q>0, \quad q^{2}-4 p r>0 \tag{16}
\end{equation*}
$$

and let $(0<) \mu_{1}<\mu_{2}$ be the roots of the equation

$$
\begin{equation*}
-p v^{2}+q v-r=0 . \tag{17}
\end{equation*}
$$

Then for all $\mu$ such that

$$
\begin{equation*}
\mu_{1}<\mu<\mu_{2}, \tag{18}
\end{equation*}
$$

the expression $-p \mu^{2}+q \mu-r$ is positive and, with respect to (14), $H_{1}>0$. Hence, the following lemma holds.

Lemma 2. Let $K_{1}, K_{3}, K_{5}$ be such that $v, \mu$ which are determined by (13) satisfy the inequalities (16) and (18) where $p, q, r$ are defined by (15) and $\mu_{1}<\mu_{2}$ are zeros of (17). Then $H_{1}$ given by (11) is positive.

Remark. A little calculation gives that the inequalities $7<v<393$ imply that $q>0$. In [7], p. 75 the implication $17<v<152 \Rightarrow q^{2}-4 p r>0$ is stated. Hence

$$
\begin{equation*}
17<v<152 \tag{19}
\end{equation*}
$$

is a sufficient condition for $v$ to satisfy (16). For such $v$ the corresponding $\mu_{1}, \mu_{2}$ are given in Fig. 4 in [7], p. 75. E.g., when $v=50$ then $\mu_{1} \doteq 0.3$ and $\mu_{2} \doteq 2.1$.

On the basis of (10), Lemma 1 implies
Theorem 4. Let $K_{1}-K_{5}$ be such that $H_{1}, H_{2}$ which are defined by (11), (12), fulfil

$$
\begin{equation*}
H_{1}-H_{2}>0 . \tag{20}
\end{equation*}
$$

Then the equilibrium point $a_{1}$ of the system (1) is conditionally stable.
Remark. When $K_{1}, K_{3}, K_{5}$ satisfy the assumption of Lemma 2 (e.g. when (19), (18) are fulfilled) and $K_{4}$ is arbitrary, then by the above argument there exists a $K_{2}^{0}>0$ such that for $K_{2}>K_{2}^{0}$ (or $K_{2}=K_{2}^{0}$ or $K_{2}<K_{2}^{0}$ ) the hypothesis (20) (or $H_{1}-H_{2}=$ $=0$ or $H_{1}-H_{2}<0$, respectively $)$ is satisfied. At the same time ( $K_{1}, K_{3}, K_{4}, K_{5}$ being fixed) there is a $K_{2}^{1}$ such that $K_{2}>K_{2}^{1}$ (or $K_{2}=K_{2}^{1}$ or $K_{2}<K_{2}^{1}$ ) implies $k<l$ (or $k=l$ or $k>l$, respectively). By Theorems 3 and $4, K_{2}^{1} \geqq K_{2}^{0}$ cannot occur. Thus $K_{2}^{1}<K_{2}^{0}$ and the following statement is true:

$$
\text { If } H_{1}-H_{2}>0, \text { then } k<l \text {. }
$$

Moreover, there is an $\varepsilon_{1}>0$ such that $k<l$ for $K_{2}^{0}-\varepsilon_{1}<K_{2}<\infty$. Since we shall consider only such $K_{2}^{\prime}$ s, we shall assume in the sequel that $k<l$. In paper [7], p. 73 even $10 k$ has been neglected with respect to $l$.
2. In order to get more information in the case when (20) is fulfilled, consider the vector field given by the system (1) in the part of the phase space which has a real meaning in the Belousov - Zhabotinskij model, i.e. for $X \geqq 0, U \geqq 0, C \geqq 0$. From the first and the third equation of (1) we conclude that
a) $\dot{X}=0$ is true in the quarterplane $X=\left(K_{5} / 5 K_{1}\right) C, 0 \leqq C<\infty, 0 \leqq U<\infty$, while $\dot{X}<0(\dot{X}>0)$ holds for $X>\left(K_{5} / 5 K_{1}\right) C,\left(X<\left(K_{5} / 5 K_{1}\right) C\right), 0 \leqq C<\infty$, $0 \leqq U<\infty($ and $0 \leqq X<\infty)$.
b) $\dot{C}=0$ is valid in the quarterplane $C=\left(K_{3} / 2 K_{5}\right) U, 0 \leqq U<\infty, 0 \leqq X<\infty$, while $\dot{C}<0(\dot{C}>0)$ holds for $C>\left(K_{3} / 2 K_{5}\right) U\left(C<\left(K_{3} / 2 K_{5}\right) U\right), 0 \leqq U<\infty$, $0 \leqq X<\infty($ and $0 \leqq C<\infty)$.

The second equation in (1) can be written in the form

$$
\dot{U}=U\left(K_{3}-K_{4} U\right)+X\left(K_{1}-K_{2} U\right)
$$

and the condition $\dot{U}=0$ leads to the quadratic equation

$$
-K_{4} U^{2}+\left(K_{3}-K_{2} X\right) U+K_{1} X=0
$$

which, for each $X \geqq 0$, possesses two real roots, a positive one $U_{p}(X)$ and a non-positive one $U_{n}(X)$. Here

$$
\begin{equation*}
U_{p}(X)=\left(K_{3}-K_{2} X+\sqrt{ }\left(\left(K_{3}-K_{2} X\right)^{2}+4 K_{1} K_{4} X\right)\right) /\left(2 K_{4}\right)(X \geqq 0) . \tag{21}
\end{equation*}
$$

Further,

$$
U_{p}^{\prime}(X)=\left(-K_{2} /\left(2 K_{4}\right)\right)\left[1+\frac{K_{3}-K_{2} X-2 K_{1} K_{4} / K_{2}}{\sqrt{ }\left(\left(K_{3}-K_{2} X\right)^{2}+4 K_{1} K_{4} X\right)}\right]
$$

Since

$$
\begin{gathered}
\left(\frac{K_{3}-K_{2} X-2 K_{1} K_{4} / K_{2}}{\sqrt{ }\left(\left(K_{3}-K_{2} X\right)^{2}+4 K_{1} K_{4} X\right)}\right)^{2}<1, \\
U_{p}^{\prime}(X)<0 \quad(X \geqq 0)
\end{gathered}
$$

while

$$
\begin{gathered}
U_{p}^{\prime \prime}(X)=\left[-2 K_{1}\left(K_{1} K_{4}-K_{3} K_{2}\right)\right] /\left[\left(K_{3}-K_{2} X\right)^{2}+4 K_{1} K_{4} X\right]^{3 / 2}>0 \\
(X \geqq 0)
\end{gathered}
$$

The inequality $U_{p}^{\prime}(X)<0$, together with $U_{p}(0)=l$, gives $l-U_{p}(X)>0$ for $0<X<\infty$. As $U_{p}$ satisfies the equality $U(l-U)\left|K_{2}+X(k-U)\right| K_{4}=0$ in which the first term is positive, we also have $k-U_{p}(X)<0(0<X<\infty)$ and thus, the inequalities

$$
\begin{equation*}
k<U_{p}(X)<l \text { for } 0<X<\infty \tag{22}
\end{equation*}
$$

are true. The results can be put together in the following statement:
c) $\dot{U}=0$ on the surface $U=U_{p}(X)$, where $U_{p}$ is determined by $(21), 0 \leqq X<\infty$,
$0 \leqq C<\infty$, while $\dot{U}<0(\dot{U}>0)$ for $U>U_{p}(X)\left(U<U_{p}(X)\right), 0 \leqq X<\infty$, $0 \leqq C<\infty$ (and $\left.0 \leqq U<\infty, X^{2}+U^{2}>0\right)$. The function $U_{p}$ is decreasing, strictly convex in $[0, \infty)$ and satisfies the inequalities (22).

The statements c), b) and a) together with (2), (3), (4) imply the following
Lemma 3. Let the numbers $U_{i}, C_{i}$ and $X_{i}, i=1,2$ be such that

$$
\begin{aligned}
& 0<U_{1} \leqq k, \quad l<U_{2}, \\
& 0<C_{1}<\left(K_{3} / 2 K_{5}\right) U_{1}, \quad\left(K_{3} / 2 K_{5}\right) U_{2}<C_{2}, \\
& 0<X_{1}<\left(K_{5} / 5 K_{1}\right) C_{1}, \quad\left(K_{5} / 5 K_{1}\right) C_{2}<X_{2} .
\end{aligned}
$$

Let $\mathscr{R}=\left\{(X, U, C): X_{1} \leqq X \leqq X_{2}, U_{1} \leqq U \leqq U_{2}, C_{1} \leqq C \leqq C_{2}\right\}$ and let $\mathscr{R}^{0}$ be its interior.

Then the following statements are true:

1. The orbit of each solution to (1) passing through a point of $\mathscr{R}$ enters $\mathscr{R}^{0}$ and remains in $\mathscr{R}^{0}$.
2. The system (1) has a unique equilibrium point in $\mathscr{R}^{0}$, namely $a_{1}$.

Lemma 3 completes the statement of Theorem 1 and implies that each solution of (1) the orbit of which passes through a point of $\mathscr{R}$ is defined on an interval $\left[t_{0}, \infty\right)$ for some $t_{0}$ (depending on that solution).

Further, for the solutions of (1) the following alternative holds:
Lemma 4. Let $(X, U, C)$ be a solution of the system (1), the orbit of which goes through a point from $\mathscr{R}$ and which is different from the equilibrium point. Then either each of its components is ultimately strictly monotone, i.e. strictly monotone in an interval $\left[t_{1}, \infty\right), t_{0}<t_{1}$, or each of its components is oscillating in the sense that its derivative changes its sign infinitely many times in each subinterval $\left[t_{1}, \infty\right)$ of $\left[t_{0}, \infty\right)$.

Proof. First, by the properties of the vector field determined by (1), it follows that if a component of the solution $(X, U, C)$ is monotone in an interval, then it is strictly monotone in that interval. Therefore it suffices to prove the following three implications:

1. If $U$ is ultimately strictly monotone (briefly u.s.m.), then $C$ is u.s.m., too.
2. If $C$ is u.s.m., then so is $X$.
3. If $X$ is u.s.m., then $U$ is also u.s.m..

Let $U$ be u.s.m., let $U$ be increasing in $\left[t_{1}, \infty\right)$. Then from the third equation of (1) it follows that $C^{\prime}(t)<0\left(C^{\prime}(t)>0\right)$ at all points $t$ where $C(t)>\left(K_{3} / 2 K_{5}\right) U(t)$ $\left(C(t)<\left(K_{3} / 2 K_{5}\right) U(t)\right)$. Since $\left(K_{3} / 2 K_{5}\right) U$ is increasing, there exist no points $t_{2}, t_{3}$ with $t_{1}<t_{2}<t_{3}$ such that $C\left(t_{i}\right)=\left(K_{3} / 2 K_{5}\right) . U\left(t_{i}\right), i=2$, 3, and $C^{\prime}(t)<0$ in $\left(t_{2}, t_{3}\right)$. Hence $C^{\prime}$ must be in a neighbourhood of $\infty$ either everywhere negative or
everywhere nonnegative. If $C$ were only nondecreasing, it would contradict the assumption that $U$ is increasing. Thus in both cases the implication 1 follows. Similarly we can proceed when $U$ is decreasing or when proving the implication 2.

The implication 3 will be proved only in the case when $X$ is increasing in $\left[t_{1}, \infty\right)$. Consider the curve $(X(t), U(t))$ in the phase plane $X, U$. By the statement c) which precedes Lemma 3, it follows that if $U(t)>U_{p}[X(t)]\left(U(t)<U_{p}[X(t)]\right)$, then $\dot{U}(t)<0(\dot{U}(t)>0)$. Since $U_{p}[X(t)]$ is decreasing, there exist no points $t_{2}, t_{3}$ with $t_{1}<t_{2}<t_{3}$ such that $U\left(t_{i}\right)=U_{p}\left[X\left(t_{i}\right)\right], i=2,3$ and $\dot{U}(t)>0$ in $\left(t_{2}, t_{3}\right)$. This gives the implication 3 in a similar way as in the previous case.

Remarks. 1. If each component of a solution $(X, U, C)$ of (1) the orbit of which goes through a point from $\mathscr{R}$ is u.s.m., then, in virtue of Lemma 3,

$$
\lim _{t \rightarrow \infty} X(t)=X_{0}, \quad \lim _{t \rightarrow \infty} U(t)=U_{0}, \quad \lim _{t \rightarrow \infty} C(t)=C_{0}
$$

2. Each of the quarterplanes mentioned in the statements $a$ ), b), as well as the surface $U=U_{p}(X)$ from the statement c), define a decomposition of the parallelepiped $\mathscr{R}$ into two subsets whereby the derivative $\dot{X}, \dot{U}, \dot{C}$ of each component of a solution to (1) has a constant in each subset of the same decomposition and these signs are mutually different in these two subsets. Hence, if each component of a solution to (1) is oscillating, then its semiorbit goes from one subset of the mentioned three decompositions into the other infinitely many times.

Lemma 5. Let $(X, U, C)$ be a solution of the system (1) the orbit of which goes through a point from $\mathscr{R}$, is different from the equilibrium point $a_{1}$ and is oscillatory in the sense of Lemma 4. Then it is oscillating around the equilibrium point $a_{1}$, i.e. the functions $x=X-X_{0}, u=U-U_{0}, c=C-C_{0}$ change their sign in each interval $\left[t_{1}, \infty\right)$ infinitely many times.

Proof. The vector function $(x, u, c)$ is a non-trivial solution of (6). It is easy to see that none of its components can be identically 0 on any interval.

Suppose now that $u(t) \geqq 0$ for $t \in\left[t_{1}, \infty\right)$. Then the third equation of the system (6) implies that for an arbitrary but fixed $t_{2}<t_{1}$ we have

$$
c(t)=c\left(t_{2}\right) \exp \left[-4 K_{5}\left(t-t_{2}\right)\right]+\int_{t_{2}}^{t} \exp \left[-4 K_{5}(t-s)\right] \cdot 2 K_{3} u(s) \mathrm{d} s
$$

This implies that either $c(t)<0$ in $\left[t_{1}, \infty\right)$ or there is a $t_{2}>t_{1}$ such that $c(t)>0$ for $t>t_{2}$.
i) If $c(t)<0$ in $\left[t_{1}, \infty\right)$, then $\dot{c}(t)>0$ in the same interval.
ii) If $c(t)>0$ in $\left(t_{2}, \infty\right)$, then similarly as in the previous case but from the first equation of (6) we conclude that either $x(t)<0$ in $\left(t_{2}, \infty\right)$ and then $\dot{x}(t)>0$ for the same $t$ or $x(t)>0$ in $\left(t_{3}, \infty\right), t_{2}<t_{3}$.
iii) The case $x(t)>0, u(t) \geqq 0$ in $\left(t_{3}, \infty\right)$, with respect to the inequalities $K_{1}$ -$-K_{2} U_{0}<0, K_{3}-K_{2} X_{0}-2 K_{4} U_{0}<0$ leads to $\dot{u}(t)<0$ in $\left(t_{3}, \infty\right)$. Hence, if $u(t) \geqq 0$ in $\left[t_{1}, \infty\right)$, then by Lemma 4 the solution $(x, u, c)$ is u.s.m. A similar implication also holds when $u(t) \leqq 0$ in $\left[t_{1}, \infty\right)$.

Therefore $u$ changes its sign in each interval $\left[t_{1}, \infty\right)$ infinitely many times.
Let now $U=U_{p}(X)$ be the function determined by (21). From (5) we get $U_{0}+$ $u=U_{p}\left(X_{0}+x\right)$, hence $u=U_{p}\left(X_{0}+x\right)-U_{0}=u_{p}(x)$. As $U_{p}\left(X_{0}\right)=U_{0}, u_{p}(0)=0$ holds, and by the statement c ) the function $u_{p}$ is decreasing and strictly convex. Consider the curve $(x, u)$. By Lemma 4, it intersects infinitely many times the graph of $u_{p}$. If all intersection points for sufficiently great $t$ lay in the halfplane $u>0$, then, again using the statement c ), we should get that $u(t)>0$ for all $t$ sufficiently great. Similarly we can exclude the case that almost all intersection points lie in the halfplane $u<0$. Therefore there exist infinitely many intersection points of the curve ( $x, u$ ) with the graph of $u_{p}$ in the halfplane $u>0$ and infinitely many of them in the halfplane $u<0$. As $u_{p}$ is decreasing and $u_{p}(0)=0$, the intersection points in the halfplane $u>0(u<0)$ also lie in the halfplane $x<0(x>0)$. So we have proved that $x$ changes its sign infinitely many times in each interval $\left[t_{1}, \infty\right)$.
If $c$ did not possess the same property, we should obtain from the first equation of (6) by a similar argument as above that $x$ is ultimately positive or ultimately negative, which gives a contradiction with the previous statement. This completes the proof of the lemma.
With help of the above lemmas we prove
Theorem 5. Let the assumption of Theorem 4 be fulfilled. Then with the exception of exactly two positive semiorbits, the positive semiorbit of each solution of (1) which is different from the equilibrium point $a_{1}$ and which goes through a point of $\mathscr{R}$ is oscillating in the sense given above and cannot tend to $a_{1}$ as $t \rightarrow \infty$. The exceptional positive semiorbits belong to the solutions of (1) each component of which is u.s.m. and that tend to $a_{1}($ as $t \rightarrow \infty)$.

Proof. Using the relation between the solutions of the systems (1) and (6) we shall prove that there exist exactly two positive semiorbits of the solutions of (6) each component of which is u.s.m.. By Remark 1, these semiorbits tend to 0 as $t \rightarrow \infty$. The positive semiorbits of all the other solutions of (6) corresponding to the solutions of (1) mentioned in the theorem are oscillatory.

Hence, consider the system (6) under the assumption of Theorem 4. By Lemma 1 it follows that the matrix

$$
A=\left(\begin{array}{ccc}
-5 K_{1}, & 0, & K_{5} \\
K_{1}-K_{2} U_{0}, K_{3}-K_{2} X_{0}-2 K_{4} U_{0}, & 0 \\
0, & 2 K_{3}, & -4 K_{5}
\end{array}\right)
$$

of the first approximation to (6) has a real eigenvalue $\alpha<0$ and a pair of complex
conjugate eigenvalues $\beta \pm i \gamma$, whereby $\beta>0$. Further, the vector function $f(z)=$ $=\left(0,-K_{2} x u-K_{4} u^{2}, 0\right)$ of the vector variable $z=(x, u, c)=\left(x_{1} \cdot x_{2}, x_{3}\right)$ satisfies the condition $\partial f / \partial z=0$ at $z=0$. Thus the system (6) satisfies all assumptions of Theorem 4.1, [1], p. 330 and, moreover, there exists a real nonsingular $3 \times 3$ constant matrix $Q$ such that (6) can be brought by the transformation $y=Q z$ to the form which already fulfils the assumptions of Theorem 8.1, [5], p. 248. By these two theorems there exists a one-dimensional manifold $S$ in the space $z$, the equations of which are ([1], p. 330)

$$
\begin{align*}
x_{i} & =q_{i 1} y_{1}+q_{i 2} \psi_{2}\left(y_{1}\right)+q_{i 3} \psi_{3}\left(y_{1}\right)  \tag{23}\\
& =x_{i}\left(y_{1}\right) \quad y_{1} \in[-\tilde{\delta}, \tilde{\delta}]
\end{align*}
$$

where $\psi_{2}, \psi_{3}$ are real analytic functions defined in $[-\tilde{\delta}, \tilde{\delta}]$,

$$
\begin{equation*}
\psi_{i}(0)=0, \quad \psi_{i}^{\prime}(0)=0 \tag{24}
\end{equation*}
$$

and the matrix $\left(q_{i k}\right)=Q^{-1}, i, k=1,2,3$, is the inverse of $Q$. With respect to (24), $x_{i}(0)=0$ and $x_{i}^{\prime}(0)=q_{i 1}, i=1,2,3$. Since $Q^{-1}$ is nonsingular, at least one of the numbers $q_{i 1}, i=1,2,3$, is different from zero. Suppose $q_{11} \neq 0$. (In the other cases we should proceed similarly.) Then we can take $\tilde{\delta}$ so small that $x_{1}\left(y_{1}\right)$ is one-to-one in $[-\tilde{\delta}, \tilde{\delta}]$.

The manifold $S$ enjoys the following property ([5], p. 248). There exist two positive numbers $\delta_{0} \leqq \delta$ such that the following statements hold:

1. If the initial point of a positive semiorbit of a solution $z=(x, u, c)$ of (6) is from $S$ and $|x| \leqq \delta_{0}$, then the whole semiorbit lies on $S$ and $\underset{t \rightarrow \infty}{\lim } z(t)=0$.
2. If a positive semiorbit of a solution $z$ of (6) lies in the $\delta$-neighbourhood of the point 0 , e.g. if $\lim _{t \rightarrow \infty} z(t)=0$, then it must lie on $S$.

Let $T$ be the set of all solutions $z$ of the system (6) each component of which is u.s.m. and for which $\lim _{t \rightarrow \infty} z(t)=0$. Let $z_{1}, z_{2} \in T$, where $z_{1}=\left(x_{1}^{*}, u_{1}, c_{1}\right), z_{2}=$ $=\left(x_{2}^{*}, u_{2}, c_{2}\right)$. Then there is an interval $\left[t_{1}, \infty\right)$ in which both functions $x_{1}^{*}, x_{2}^{*}$ are one-to-one and possess a constant sign. Let both functions $x_{1}^{*}, x_{2}^{*}$ be positive in this interval. If we denote by $t_{i}$ the inverse function to $x_{i}^{*}, i=1,2, t_{i}$ maps the interval $\left(0, x_{i}^{*}\left(t_{1}\right)\right]$ onto $\left[t_{1}, \infty\right)$ and a positive semiorbit corresponding to the solution $z_{i}$ is given by the relations

$$
\begin{equation*}
u=u_{i}\left[t_{i}(x)\right], \quad c=c_{i}\left[t_{i}(x)\right], \quad x \in\left(0, x_{i}^{*}\left(t_{1}\right)\right], \quad i=1,2 . \tag{25}
\end{equation*}
$$

On the other hand, since $\lim _{t \rightarrow \infty} z_{i}(t)=0$, this semiorbit must lie on $S$ (taking $t_{1}$ sufficiently large). Hence, in virtue of (23) we have

$$
\begin{equation*}
x=x_{1}\left(y_{1}\right), \quad u=x_{2}\left(y_{1}\right), \quad c=x_{3}\left(y_{1}\right), \quad y_{1} \in[-\tilde{\delta}, \tilde{\delta}] . \tag{26}
\end{equation*}
$$

Now we choose an $x \in\left(0, x_{1}^{*}\left(t_{1}\right)\right] \cap\left(0, x_{2}^{*}\left(t_{1}\right)\right]$. As $x_{1}$ is injective, there exists exactly one $y_{1}$ such that $x=x_{1}\left(y_{1}\right)$. Comparing (25) with (26) we get at the chosen point

$$
u_{1}\left[t_{1}(x)\right]=x_{2}\left(y_{1}\right)=u_{2}\left[t_{2}(x)\right], \quad c_{1}\left[t_{1}(x)\right]=x_{3}\left(y_{1}\right)=c_{2}\left[t_{2}(x)\right]
$$

and hence both semiorbits pass through the same point which implies that they are identical (for all $\left.x \in\left(0, x_{1}^{*}\left(t_{1}\right)\right] \cap\left(0, x_{2}^{*}\left(t_{1}\right)\right]\right)$. The set $T$ can have at most two semiorbits (one for $x>0$ and another for $x<0$ ).

Further, we prove that no oscillatory solution $z$ of (6) can satisfy $\lim _{t \rightarrow \infty} z(t)=0$. Otherwise a positive semiorbit of this solution would lie on $S$ and since $z$ is oscillating, two points $t_{1}, t_{2}$ would exist with $t_{1}<t_{2}$, such that $z\left(t_{1}\right)=z\left(t_{2}\right)$. This would imply that $z$ is periodic which contradicts the fact that $\lim _{t \rightarrow \infty} z(t)=0$. There is even no positive semiorbit of an oscillatory solution $z$ of (6) lying entirely in a $\delta$-neighbourhood of 0 where $\delta>0$ is sufficiently small. Hence only positive semiorbits of the solutions of (6) each component of which is u.s.m. may lie on $S$. From the property 1 of $S$ it follows that such semiorbits do exist. This completes the proof of Theorem 5.

Further properties of oscillatory solutions of (1) are given by
Theorem 6. Let the assumptions of Theorem 4 be fulfilled. Le ( $X, U, C$ ) be a solution of the system (1) the orbit of which goes through a point from $\mathscr{R}$, is different from the equilibrium point $a_{1}$ and is oscillatory in the sense of Lemma 4 and Lemma 5, respectively. Then the following statement holds for the functions

$$
x=X-X_{0}, \quad u=U-U_{0}, \quad c=C-C_{0}
$$

in each interval $\left(t_{1}, \infty\right)$ where $t_{1}$ is a zero-point of one of them: Between two successive zeros of $x$ and $u$ and $c$, there exists exactly one zero of $u$ and $c$ and $x$ respective$l y$. All zeros of the functions $x, u, c$ in $\left(t_{1}, \infty\right)$ are simple.

Proof. The function $(x, u, c)$ is a solution of (6) which is defined in a neighbourhood of $\infty$, say in ( $t_{0}, \infty$ ) where $t_{0}<t_{1}$. Let $t_{2}, t_{3} \in R$ be such that $t_{1} \leqq t_{3} \leqq t_{2}$, $t_{1}<t_{2}$. Directly from (6) we get the statement:

1. None of the functions $x, u, c$ possesses a zero of multiplicity 3 (and hence all zeros of these functions are isolated and cannot have a limit point in $\left(t_{0}, \infty\right)$ ), no pair of these functions possess a double zero at the same point. Further, $u\left(t_{2}\right)=$ $=\dot{u}\left(t_{2}\right)=0$ implies $x\left(t_{2}\right)=0, x\left(t_{2}\right)=\dot{x}\left(t_{2}\right)=0$ gives $c\left(t_{2}\right)=0$ and $c\left(t_{2}\right)=$ $=\dot{c}\left(t_{2}\right)=0$ yields the equality $u\left(t_{2}\right)=0$.
The following relations between $x, u, c$ can be derived from (6):

$$
\begin{align*}
& x(t)=x\left(t_{3}\right) \exp \left[-5 K_{1}\left(t-t_{3}\right)\right]+\int_{t_{3}}^{t} \exp \left[-5 K_{1}(t-s)\right] K_{5} c(s) \mathrm{d} s  \tag{27}\\
& c(t)=c\left(t_{3}\right) \exp \left[-4 K_{5}\left(t-t_{3}\right)\right]+\int_{t_{3}}^{t} \exp \left[-4 K_{5}(t-s)\right] 2 K_{3} u(s) \mathrm{d} s
\end{align*}
$$

$$
\begin{align*}
& u(t)=u\left(t_{3}\right) \exp \left[\left(K_{3}-K_{2} X_{0}-2 K_{4} U_{0}\right)\left(t-t_{3}\right)-K_{2} \int_{t_{3}}^{t} x(s) \mathrm{d} s\right]+  \tag{29}\\
& \quad+\int_{t_{3}}^{t} K(t, s)\left[\left(K_{1}-K_{2} U_{0}\right) x(s)-K_{4} u^{2}(s)\right] \mathrm{d} s, \quad t_{0}<t<\infty
\end{align*}
$$

where $K=K(t, s)$ is the Cauchy function of the differential equation $\dot{u}=\left[K_{3}-\right.$ $\left.-K_{2} X_{0}-2 K_{4} U_{0}-K_{2} x(t)\right] u$ and, hence, positive.
Owing to the inequality $K_{1}-K_{2} U_{0}<0$ as well as to the fact that the zeros of $x$, $u, c$ are isolated, we conclude from (27)-(29):
2. If $u(t) \geqq 0(u(t) \leqq 0)$ in $\left[t_{1}, t_{2}\right]$, then either $c$ does not posses any zero in $\left[t_{1}, t_{2}\right]$ or there is a point $t_{3} \in\left[t_{1}, t_{2}\right]$ such that $c\left(t_{3}\right)=0$ and, if $t_{3}<t_{2}$, then $c(t)>0$ $(c(t)<0)$ in $\left(t_{3}, t_{2}\right]$ and, if $t_{1}<t_{3}$, then $c(t)<0(c(t)>0)$ in $\left[t_{1}, t_{3}\right)$.
3. If $c(t) \geqq 0(c(t) \leqq 0)$ in $\left[t_{1}, t_{2}\right]$, then either $x$ does not possess any zero in $\left[t_{1}, t_{2}\right]$, or there is a point $t_{3} \in\left[t_{1}, t_{2}\right]$ such that $x\left(t_{3}\right)=0$ and, if $t_{3}<t_{2}$, then $x(t)>$ $>0(x(t)<0)$ in $\left(t_{3}, t_{2}\right]$ and, if $t_{1}<t_{3}$, then $x(t)<0(x(t)>0)$ in $\left[t_{1}, t_{3}\right)$.
4. If $x(t) \geqq 0$ in $\left[t_{1}, t_{2}\right]$, then either $u$ does not have any zero in $\left[t_{1}, t_{2}\right]$, or there is a point $t_{3} \in\left[t_{1}, t_{2}\right]$ such that $u\left(t_{3}\right)=0$ and, if $t_{3}<t_{2}$, then $u(t)<0$ in $\left(t_{3}, t_{2}\right]$ and, if $t_{1}<t_{3}$, then $u(t)>0$ in $\left[t_{1}, t_{3}\right)$.

The next statement is based on the properties of the vector field (6).
5. If $x(t) \leqq 0$ in $\left[t_{1}, t_{2}\right]$, then either $u$ does not possess any zero in $\left[t_{1}, t_{2}\right]$ or there is a point $t_{3} \in\left[t_{1}, t_{2}\right]$ such that $u\left(t_{3}\right)=0$ and, if $t_{3}<t_{2}$, then $u(t)<0$ in $\left(t_{3}, t_{2}\right]$ and, if $t_{1}<t_{3}$, then $u(t)>0$ in $\left[t_{1}, t_{3}\right)$.

Proof. Consider the curve $(x, u)$ in the phase plane $x, u$. If $u_{p}$ has the same meaning as in the proof of Lemma 5, then by the statement c) preceding Lemma 3, the function $u$ is decreasing (increasing) in the intervals for which the points ( $x, u$ ) lie above (below) the graph of $u_{p}$.

Let $t_{3}$ be the first zero of $u$ in $\left[t_{1}, t_{2}\right]$. If $x\left(t_{3}\right)=0$ then by (6), $\dot{u}\left(t_{3}\right)=0$. As $t_{3}$ cannot be a double zero of both functions $x, u$, either $\dot{x}\left(t_{3}\right)>0$ and then $t_{3}=t_{2}$ or $\dot{x}\left(t_{3}\right)<0$ and then $t_{3}=t_{1}$, with respect to $x(t) \leqq 0$ in $\left[t_{1}, t_{2}\right]$. In the first case $\left(\dot{x}\left(t_{3}\right)>0, x\left(t_{3}\right)=u\left(t_{3}\right)=\dot{u}\left(t_{3}\right)=0\right)$, the properties of the vector field imply that $\ddot{u}\left(t_{3}\right)>0$ cannot happen. Hence $\ddot{u}\left(t_{3}\right)<0$ and it follows again by these properties that $u(t)<0$ must hold in $\left[t_{1}, t_{3}\right)$. In the second case $\ddot{u}\left(t_{3}\right)<0$ cannot happen, therefore $\ddot{u}\left(t_{3}\right)>0$ and hence $u(t)>0$, first in a right neighbourhood of $t_{3}$ and then, by using the properties of the vector field (6) as well as the result from the first case, we come to the conclusion that $u(t)>0$ in $\left(t_{3}, t_{2}\right]$.

If $x\left(t_{3}\right)<0$, then we conclude again by the properties of the vector field (6) that $u(t)<0$ in $\left[t_{1}, t_{3}\right)$ and by the above result from the first case, $u(t)>0$ in $\left(t_{3}, t_{2}\right]$. This completes the proof of the statement 5 .

Without assuming that $t_{1}$ is a zero of one of the functions $x, u, c$ we shall prove the statement
6. If either $x$ or $u$ or $c$, respectively, has a double zero at $t_{2}$, then the functions $x, u, c$ do not vanish at any $t<t_{2}$ and $x(t) u(t)>0, x(t) c(t)<0$ for all $t$ such that $t_{0}<t<t_{2}$.

Proof. Let $u\left(t_{2}\right)=\dot{u}\left(t_{2}\right)=0$. Then, by the statement $1, x\left(t_{2}\right)=0, c\left(t_{2}\right) \dot{x}\left(t_{2}\right) \ddot{u}\left(t_{2}\right) \neq$ $\neq 0$. Statements 4 and 5 imply

$$
\begin{equation*}
\dot{x}\left(t_{2}\right) \ddot{u}\left(t_{2}\right)<0 \tag{30}
\end{equation*}
$$

(here the behaviour of the functions $x$ and $u$ to the right from $t_{2}$ is considered). The statement 3 implies

$$
\begin{equation*}
\dot{x}\left(t_{2}\right) c\left(t_{2}\right)>0 . \tag{31}
\end{equation*}
$$

Suppose that $u$ vanishes to the left from $t_{2}$, hence there exists a $t_{1}<t_{2}$ such that $u(t) \neq 0$ in $\left(t_{1}, t_{2}\right)$ and $u\left(t_{1}\right)=0$. As $x\left(t_{2}\right)=u\left(t_{2}\right)=0$, by the statements 4 and 5 we have that there is a $t_{3}$ such that $t_{1}<t_{3}<t_{2}$ with $x(t) \neq 0$ in $\left(t_{3}, t_{2}\right)$ and $x\left(t_{3}\right)=0$. Because $x\left(t_{2}\right)=0$ and $c\left(t_{2}\right) \neq 0$, by the statement 3 there is a $t_{4}$ with $t_{3}<t_{4}<t_{2}$ such that $c(t) \neq 0$ in $\left(t_{4}, t_{2}\right)$ and $c\left(t_{4}\right)=0$. This contradicts the statement 2 . Similarly we come to a contradiction when we suppose the existence of a point $t_{3}$ or $t_{4}$, respectively, with the above properties (i.e. $x\left(t_{3}\right)=0, x(t) \neq 0$ in $\left(t_{3}, t_{2}\right), c\left(t_{4}\right)=0$ and $c(t) \neq 0$ in $\left.\left(t_{4}, t_{2}\right)\right)$. Hence $x(t) u(t) c(t) \neq 0$ for all $t, t_{0}<t<t_{2}$ and on the basis of (30), (31) we can state more precisely that $x(t) u(t)>0, x(t) c(t)<0$ for $t_{0}<t<$ $<t_{2}$.

If $x\left(t_{2}\right)=\dot{x}\left(t_{2}\right)=0$, then $c\left(t_{2}\right)=0$ by the statement 1 and using statements 3 and 2 we get $\dot{c}\left(t_{2}\right) \ddot{x}\left(t_{2}\right)>0$ and $\dot{c}\left(t_{2}\right) u\left(t_{2}\right)>0$.

Let there exist a point $t_{1}<t_{2}$ such that $x\left(t_{1}\right)=0, x(t) \neq 0$ in $\left(t_{1}, t_{2}\right)$. As $x\left(t_{2}\right)=$ $=c\left(t_{2}\right)=0$, by 3 there is a $t_{3}, t_{1}<t_{3}<t_{2}$ such that $c\left(t_{3}\right)=0, c(t) \neq 0$ in $\left(t_{3}, t_{2}\right)$. Because $c\left(t_{2}\right)=0$ and $u\left(t_{2}\right) \neq 0$, by the statement 2 there is a $t_{4}, t_{3}<t_{4}<t_{2}$ such that $u\left(t_{4}\right)=0, u(t) \neq 0$ in $\left(t_{4}, t_{2}\right)$. This contradicts the statement 4. Hence the statement 6 is true in this case.

When $c\left(t_{2}\right)=\dot{c}\left(t_{2}\right)=0$, we have $\ddot{c}\left(t_{2}\right) \dot{u}\left(t_{2}\right)>0, x\left(t_{2}\right) . \dot{u}\left(t_{2}\right)<0$ in virtue of the statements 2, 4 and 5. If there is a $t_{1}<t_{2}$ such that $c\left(t_{1}\right)=0, c(t) \neq 0$ in $\left(t_{1}, t_{2}\right)$, then there are $t_{1}<t_{3}<t_{4}<t_{2}$ with $u\left(t_{3}\right)=0, x\left(t_{4}\right)=0, u(t) \neq 0$ in $\left(t_{3}, t_{2}\right)$, $x(t) \neq 0$ in $\left(t_{4}, t_{2}\right)$. This leads to contradiction with the statement 3. Again 6 is true.
7. Between two successive zeros of $x$ and $u$ and $c$, there exists in $\left(t_{1}, \infty\right)$ exactly one zero of $u$ and $c$ and $x$, respectively, and all zeros of the functions $x, u, c$ are simple.

Only the first assertion will be proved. The other two can be proved in a similar way. Since $t_{1}$ is a zero of one of the functions $x, u, c$, the statement 6 implies that all zeros of $x, u, c$ in $\left(t_{1}, \infty\right)$ are simple and hence (6) implies that none two of the functions $x, u, c$ have a common zero-point.

Let $t_{1}<s_{1}<s_{2}, x\left(s_{1}\right)=x\left(s_{2}\right)=0$ and $x(t) \neq 0$ in $\left(s_{1}, s_{2}\right)$. By the statements 4 and $5, u$ can have at most one zero in ( $s_{1}, s_{2}$ ). Suppose it has none. Then the same
is true for $\left[s_{1}, s_{2}\right]$ and two cases will be distinguished:
i) $\operatorname{sgn} x(t)=\operatorname{sgn} u(t) \quad\left(t \in\left(s_{1}, s_{2}\right)\right)$.

Let $s_{3}>s_{2}$ be the next zero of $x$, i.e. $x\left(s_{3}\right)=0, x(t) \neq 0$ in $\left(s_{2}, s_{3}\right)$. By the statements 4 and $5, u(t) \neq 0$ in $\left[s_{1}, s_{3}\right]$. Then 2 gives that $c$ possesses at most one zero in $\left[s_{1}, s_{3}\right]$, either in $\left(s_{1}, s_{2}\right)$ or in $\left(s_{2}, s_{3}\right)$. In all three cases (the first case occurs when $c$ does not vanish in $\left[s_{1}, s_{3}\right]$ ) we come with regard to the statement 3 to contradiction.
ii) $\operatorname{sgn} x(t) \neq \operatorname{sgn} u(t) \quad\left(t \in\left(s_{1}, s_{2}\right)\right)$.

Here two subcases may occur. If $x$ has no zero smaller than $s_{1}$ and greater or equal to $t_{1}$, then by statements 4 and 5 , neither $u$ has. By $2, c$ possesses at most one zero in $\left[t_{1}, s_{2}\right]$ and, by the meaning of $t_{1}, t_{1}$ is the unique zero of $c$ in $\left[t_{1}, s_{2}\right]$. But $x\left(s_{1}\right)=$ $=x\left(s_{2}\right)=0$ leads to a contradiction with the statement 3 .
If $s_{0}<s_{1}$ is the greatest zero of $x$ smaller than $s_{1}$, then by the statements 4 and 5 , $u(t) \neq 0$ in $\left[s_{0}, s_{2}\right]$. This leads to a contradiction similarly as in the case i).

The contradiction shows that between $s_{1}$ and $s_{2}$ there exists a unique zero of $u$.
3. Finally we shall prove existence of a periodic solution to (1). To that aim we shall apply Hopf's theorem and the method from the paper [4].

Theorem 7. Let $K_{1}, K_{3}, K_{5}$ be such that $H_{1}$ determined by (11) is positive. Let $K_{4}$ be arbitrary and $K_{2}=K_{2}^{0}$ be such that $H_{1}-H_{2}=0$, where $H_{2}$ is defined by (12). Then there exists an $\varepsilon_{0}>0$ such that for $K_{2} \in\left(K_{2}^{0}-\varepsilon_{0}, K_{2}^{0}+\varepsilon_{0}\right)$ the system (1) has a periodic solution different from the equilibrium point whose orbit lies in $R$ exactly in one of the three cases: Either only for each $K_{2} \in\left(K_{2}^{0}-\varepsilon_{0}, K_{2}^{0}\right)$ or only for each $K_{2} \in\left(K_{2}^{0}, K_{2}^{0}+\varepsilon_{0}\right)$ or only for $K_{2}=K_{2}^{0}$.

Proof. Put $K_{2}=K_{2}^{0}+\varepsilon$. As the coefficients of (7) determined by $\left(8_{1}\right)$ are continuous in $K_{2}$, the eigenvalues of the matrix of the system of the first approximation are continuous in $K_{2}$. Hence, by Lemma 1, there is an $\varepsilon_{0}>0$ such that for $|\varepsilon|<\varepsilon_{0}$ one eigenvalue $-\bar{a}(\varepsilon)$ is real and negative, while the other two are complex conjugate. Denote them $e(\varepsilon) \pm i(\bar{b}(\varepsilon))^{1 / 2}$, where $\bar{a}(0)=a(0), \bar{b}(0)=b(0)$ and $a(0), b(0)$ stand for the coefficients $a, b$ in $\left(8_{1}\right)$ for $\varepsilon=0$.

Differentiating (7) with respect to $\varepsilon$, we get

$$
\begin{gathered}
\lambda^{\prime}(\varepsilon)=\left[-a^{\prime}(\varepsilon) \lambda^{2}(\varepsilon)-b^{\prime}(\varepsilon) \lambda(\varepsilon)-d^{\prime}(\varepsilon)\right] /\left[3 \lambda^{2}(\varepsilon)+\right. \\
+2 a(\varepsilon) \lambda(\varepsilon)+b(\varepsilon)] .
\end{gathered}
$$

Since for $\varepsilon=0$ the roots of (7) are distinct, the Implicit Function Theorem guarantees the existence of $\lambda^{\prime}(\varepsilon)$ for sufficiently small $\varepsilon$. But $\lambda^{\prime}(\varepsilon)=e^{\prime}(\varepsilon) \pm(1 / 2)[\bar{b}(\varepsilon)]^{-1 / 2} \bar{b}^{\prime}(\varepsilon) i$, therefore

$$
\begin{gathered}
e^{\prime}(0)=\operatorname{Re}\left\{\left[a^{\prime}(0) b(0)-d^{\prime}(0)-b^{\prime}(0) .\right.\right. \\
\left.\left.\cdot\left( \pm i[b(0)]^{1 / 2}\right)\right] /\left[-2 b(0)+2 a(0)\left( \pm i[b(0)]^{1 / 2}\right)\right]\right\}=
\end{gathered}
$$

$$
\begin{gathered}
=\left(2 b(0) /\left[4 b^{2}(0)+4 a^{2}(0) b(0)\right]\right)\left[-a^{\prime}(0) b(0)+\right. \\
\left.d^{\prime}(0)-b^{\prime}(0) a(0)\right] .
\end{gathered}
$$

Here $\left(2 b(0) /\left[4 b^{2}(0)+4 a^{2}(0) b(0)\right]\right)>0, a^{\prime}(0)=-11 K_{1} K_{3}^{2} K_{4} /\left(K_{2}^{0} K_{3}+10 K_{1} K_{4}\right)^{2}$, $d^{\prime}(0)=0$ and $b^{\prime}(0)=-K_{1} K_{3}^{2} K_{4}\left(55 K_{1}+44 K_{5}\right) /\left(K_{2}^{0} K_{3}+10 K_{1} K_{4}\right)$ and thus $e^{\prime}(0)>$ $>0$. All hypotheses of Hopf's theorem being satisfied, by this theorem there exist periodic solutions of (1) with orbits in a neighbourhood of the equilibrium point $a_{1}$ for $K_{2}=K_{2}^{0}$ exactly in one of the three cases: Either for $\varepsilon>0$, or for $\varepsilon=0$, or for ह $<0$.

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## Súhrn

## O EXISTENCII OSCILATORICKÝCH RIEŠENÍ VO WEISBUCHOVOM-SALOMONOVOM-ATLANOVOM MODELE BELOUSOVEJ-ŽABOTINSKÉHO REAKCIE

Valter Šeda

V práci sa pojednáva o stabilite riešení diferenciálneho systému (1), ktorý predstavuje uvažovaný model Belousovej-Žabotinského reakcie. Dokazuje sa existencia oscilatorických riešení tohto systému a veta o oddelovaní nulových bodov jednotlivých zložiek takýchto riešení. Záverom je dokázaná aj existencia periodického riešenia.

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