# Asymptotic behavior and stability of second order neutral delay differential equations 

Guiling Chen<br>Department of Mathematics, Southwest Jiaotong University, Chengdu 610031, PR China<br>Onno van Gaans<br>Mathematical Institute, Leiden University, P.O.Box 9512, 2300 RA, Leiden, The Netherlands<br>Sjoerd Verduyn Lunel<br>Mathematical Institute, Utrecht University, PO Box 80010, 3508 TA Utrecht, The Netherlands


#### Abstract

We study the asymptotic behavior of a class of second order neutral delay differential equations by both a spectral projection method and an ordinary differential equation method approach. We discuss the relation of these two methods and illustrate some features using examples. Furthermore, a fixed point method is introduced as a third approach to study the asymptotic behavior. We conclude the paper with an application to a mechanical model of turning processes.


Keywords: Asymptotic behavior, stability, asymptotic stability, neutral delay differential equation, retarded delay differential equation, characteristic equation, fixed point theory, spectral theory.

## 1. Introduction

Neutral delay differential equations arise from a variety of applications including control systems, electrodynamics, mixing liquids, neutron transportation and population models. In the qualitative analysis of such systems, the stability and asymptotic behavior of solutions play an important role. In 1973, Driver, Sasser and Slater [4] studied asymptotic behavior, oscillation and stability of first order delay differential equations with small delay using an approach based on an ordinary differential equation (ODE) method. The key idea of the ODE approach is to transform the differential equation into a lower order equation by using a real root of the corresponding characteristic equation. Following this approach as presented in [4], a number of papers appeared in which the asymptotic behavior, oscillation and stability for first (or second or higher) order (neutral) delay differential equations, and integro-differential equations with unbounded delay as well as for delay difference equations were studied, see $[9,13,14,15,12]$. A disadvantage of this ODE approach is that it does not lead to explicit formulas for the reduced lower order equations.

In 2003, by using spectral theory, Frasson and Verduyn Lunel [5] presented a new approach to study the asymptotic behavior of neutral delay differential equations, the so-called spectral projection method.

In this paper it is our intention to compare the two approaches. We discuss their relations by studying asymptotic behavior of a class of second order neutral delay differential equations. We obtain that under the same assumptions, the ODE approach is equivalent to the spectral approach (see Section 4). However, the spectral approach has some advantages, since the conditions for the spectral method are weaker than those needed for the ODE method, as is illustrated by Example 4.7, and the asymptotic behavior of neutral delay differential equations can be presented by a general formula (see Theorem 2.5). Furthermore, by using the spectral approach, we can also study the asymptotic behavior of neutral delay differential equations with matrix coefficients.

[^0]We consider a specific class of second order neutral delay differential equations of the following form

$$
\left\{\begin{array}{l}
x^{\prime \prime}(t)+c x^{\prime \prime}(t-\tau)=p_{1} x^{\prime}(t)+p_{2} x^{\prime}(t-\tau)+q_{1} x(t)+q_{2} x(t-\tau),  \tag{1}\\
x(t)=\phi(t), \quad-\tau \leq t \leq 0
\end{array}\right.
$$

where $c, p_{1}, p_{2}, q_{1}, q_{2} \in \mathbb{R}, \tau>0$, the initial function $\phi$ is a given continuously differentiable real-valued function on the initial interval $[-\tau, 0]$.

Such neutral delay differential equations arise, for example, in the study of interconnected oscillatory systems described by a system of linear hyperbolic partial differential equations. Here the motion of each of the individual oscillatory systems is described by a boundary condition and the interconnection is given by a travelling wave. Generally, such systems can be transformed into delay equations involving delays in the highest derivative, see [9]. Furthermore, neutral delay differential equations serve as population models in the sense that they can be interpreted as special cases of the standard Gurtin-MacCamy model for a population structured by age with birth and death rate depending on the total adult population, see [8].

Equation (1) is a general second order neutral differential equation with constant coefficients, where all delays in the delayed terms are assumed to be equal and the coefficient in front of $x^{\prime \prime}(t)$ is assumed to be nonzero and scaled to 1. The second order equation is rich enough to illustrate the analysis of equations of order higher than one and simple enough to allow for a clear comparison of different approaches.

A special case of system (1) is a retarded delay equation, i.e.,

$$
\begin{equation*}
x^{\prime \prime}(t)+a x^{\prime}(t)+b x(t-r)+c x(t)=0, \quad a, b, c \in \mathbb{R}, \quad r>0 \tag{2}
\end{equation*}
$$

which is often called a delayed oscillator, is well-studied in applications. It appears, for example, as the basic governing equation of the regenerative model of machine tool chatter. We illustrate a third approach, based on a fixed point method, to study the asymptotic behavior of such equations.

The organization of this paper is as follows. In Section 2, the spectral approach is introduced and used to study the asymptotic behavior of the solutions of (1). In Section 3, the ODE approach is introduced to study the asymptotic behavior of solutions of (1). In Section 4, both approaches are analysed by investigating a number of examples. Finally, in Section 5, we present an approach based on the fixed point method and use this approach to study the asymptotic behavior of (2). As an application, a mechanical model of turning processes is presented in Section 6.

## 2. Asymptotic behavior by spectral approach

Let $C=C\left([-\tau, 0], \mathbb{C}^{n}\right)$ denote the Banach space of continuous functions endowed with the supremum norm. From the Riesz representation theorem it follows that every bounded linear mapping $L: C \rightarrow \mathbb{C}^{n}$ can be represented by

$$
L \varphi=\int_{-\tau}^{0} d \eta(\theta) \varphi(\theta)
$$

where $\eta(\theta),-\tau \leq \theta \leq 0$, is an $n \times n$-matrix whose elements are of bounded variation, normalized so that $\eta$ is continuous from the left on $(-\tau, 0)$ and $\eta(0)=0$, shortly, $\eta \in N B V$. For a function $x:[-\tau, \infty) \rightarrow \mathbb{C}^{n}$, we denote by $x_{t} \in C$ the function $x_{t}(\theta)=x(t+\theta),-\tau \leq \theta \leq 0$ and $t \geq 0$.

An initial value problem for a linear autonomous neutral functional differential equation (NFDE) is given by the following relation

$$
\left\{\begin{array}{l}
\frac{d}{d t} D x_{t}=L x_{t}, \quad t \geq 0,  \tag{3}\\
x_{0}=\phi, \quad \phi \in C,
\end{array}\right.
$$

where $D: C \rightarrow \mathbb{C}^{n}$ is continuous, linear and atomic at zero, $L: C \rightarrow \mathbb{C}^{n}$ is linear and continuous and, both operators are respectively, presented by

$$
L \varphi=\int_{-\tau}^{0} d \eta(\theta) \varphi(\theta), \quad D \varphi=\varphi(0)-\int_{-\tau}^{0} d \mu(\theta) \varphi(\theta)
$$

where $\eta, \mu \in N B V\left([-\tau, 0], \mathbb{C}^{n \times n}\right)$, and $\mu$ is continuous at zero. See Hale and Verduyn Lunel [9] for a detailed information.

For the second order neutral functional differential equation (1), let $y(t)=x^{\prime}(t)$, then (1) can be written in the form

$$
\left\{\begin{array}{l}
x^{\prime}(t)=y(t) \\
y^{\prime}(t)+c y^{\prime}(t-\tau)=p_{1} y(t)+p_{2} y(t-\tau)+q_{1} x(t)+q_{2} x(t-\tau) .
\end{array}\right.
$$

If $X(t)=\binom{x(t)}{y(t)}$, then we have

$$
\begin{equation*}
X^{\prime}(t)+C X^{\prime}(t-\tau)=E X(t)+F X(t-\tau) \tag{4}
\end{equation*}
$$

where

$$
C=\left(\begin{array}{ll}
0 & 0 \\
0 & c
\end{array}\right), E=\left(\begin{array}{cc}
0 & 1 \\
q_{1} & p_{1}
\end{array}\right) \quad \text { and } \quad F=\left(\begin{array}{cc}
0 & 0 \\
q_{2} & p_{2}
\end{array}\right) .
$$

By taking $\mu(\theta)=C$, for $\theta \leq-\tau, \mu(\theta)=0$, for $\theta>-\tau$, and $\eta(\theta)=-F$, for $\theta \leq-\tau, \eta(\theta)=0$, for $-\tau<\theta<0, \eta(\theta)=E$, for $\theta \geq 0$, (1) can be written in the form (3).

Throughout this paper, a continuous real-valued function $x$ defined on the interval $[-\tau, \infty)$ is said to be a solution of the initial value problem (1) if $x$ satisfies (1) in the mild sense, see Lemma 2.1. It is well known (see [4]) that for any given initial function $\phi$, there exists a unique solution of the initial value problem (1).

Given the solution $x(\phi)$ of the initial value problem (3), the solution operator $T(t): C \rightarrow C$ is defined by the relation

$$
T(t) \phi=x_{t}(. ; \phi), \quad t \geq 0
$$

Lemma 2.1. (Hale and Verduyn Lunel [9]) The solution operator $T(t)$ is a $C_{0}$-semigroup on $C$ with infinitesimal generator

$$
\left\{\begin{array}{l}
D(A)=\left\{\phi \in C \left\lvert\, \frac{d \phi}{d \theta} \in C\right., D \frac{d \phi}{d \theta}=L \phi\right\}  \tag{5}\\
A \phi=\frac{d \phi}{d \theta}
\end{array}\right.
$$

Lemma 2.2. (Hale and Verduyn Lunel [9]) If $A$ is defined by equation (5), then $\sigma(A)=P_{\sigma(A)}$ and $\lambda \in \sigma(A)$ if and only if $\lambda$ satisfies the characteristic equation $\operatorname{det} \Delta(\lambda)=0$, where

$$
\begin{equation*}
\Delta(\lambda)=\lambda I-\int_{-\tau}^{0} \lambda e^{\lambda \theta} d \mu(\theta)-\int_{-\tau}^{0} e^{\lambda \theta} d \eta(\theta) \tag{6}
\end{equation*}
$$

where $P_{\sigma(A)}$ denotes the point spectrum of $A$.
It is well known that there is a close connection between the spectral properties of the inifinitesimal generator $A$ and the characteristic matrix $\Delta(\lambda)$ given by (6). In particular, the geometric multiplicity $d_{\lambda}$ is equal to the dimension of the null space of $\Delta(z)$ at $z=\lambda$, and the algebraic multiplicity $m_{\lambda}$ is equal to the multiplicity of $z=\lambda$ as a zero of $\operatorname{det} \Delta(\lambda)=0$. Furthermore, the generalized eigenspace at $\lambda$ is given by

$$
\mathcal{M}_{\lambda}=\mathcal{N}(\lambda I-A)^{k_{\lambda}},
$$

where $k_{\lambda}$ denotes the order of $z=\lambda$ as a pole of $\Delta(z)^{-1}$. See Lemma 2.1 on page 263 of [9].
Lemma 2.3. (Hale and Verduyn Lunel [9]) For any $\lambda$ in $\sigma(A)$, the generalized eigenspace $\mathcal{M}_{\lambda}(A)$ is finite dimensional and there is an integer $k$ such that $\mathcal{M}_{\lambda}(A)=\mathcal{N}\left((\lambda I-A)^{k}\right)$ and we have a direct sum decomposition

$$
C=\mathcal{N}\left((\lambda I-A)^{k}\right) \oplus \mathcal{R}\left((\lambda I-A)^{k}\right)
$$

From the spectral theory [9,3], it follows that the spectral projection onto $\mathcal{M}_{\lambda}(A)$ along $\mathcal{R}\left((\lambda I-A)^{k}\right)$ can be represented by a Dunford integral

$$
\begin{equation*}
P_{\lambda}=\frac{1}{2 \pi i} \int_{\Gamma_{\lambda}}(z I-A)^{-1} d z \tag{7}
\end{equation*}
$$

where $\Gamma_{\lambda}$ is a small circle such that $\lambda$ is the only singularity of $(z I-A)^{-1}$ inside $\Gamma_{\lambda}$. In the following, the main results of the explicit representation of asymptotic bahavior of neutral functional differential equations in Frasson and Verduyn Lunel [5] is introduced.

Definition 2.4. An eigenvalue $\lambda_{d}$ is called a dominant eigenvalue of $A$, if there exists $a \epsilon>0$, such that if $\lambda$ is another eigenvalue of $A$, then $\operatorname{Re} \lambda<\operatorname{Re} \lambda_{d}-\epsilon$.

Theorem 2.5. (Frasson and Verduyn Lunel [5]) Let A be given by (5), if A has a simple and dominant eigenvalue $\lambda_{d}$, then there exists positive numbers $\epsilon$ and $M$ such that

$$
\left\|e^{-\lambda_{d} t} T(t) \phi-P_{\lambda_{d}} \phi\right\| \leq M e^{-\epsilon t}
$$

and

$$
\lim _{t \rightarrow \infty} e^{-\lambda_{d} t} T(t) \phi=e^{\lambda_{d}} \cdot\left[\frac{d}{d z} \operatorname{det} \Delta\left(\lambda_{d}\right)\right]^{-1} \operatorname{adj} \Delta\left(\lambda_{d}\right) K\left(\lambda_{d}\right) \phi
$$

Furthermore, if $x(t)=x(\cdot, \phi)$ denotes the solution of (3) with initial data $x_{0}=\phi$, then

$$
\lim _{t \rightarrow \infty} e^{-\lambda_{d} t} x(t)=\left[\frac{d}{d z} \operatorname{det} \Delta\left(\lambda_{d}\right)\right]^{-1} \operatorname{adj} \Delta\left(\lambda_{d}\right) K\left(\lambda_{d}\right) \phi
$$

where $\operatorname{adj} \Delta\left(\lambda_{d}\right)$ denotes the matrix of cofactors of $\Delta\left(\lambda_{d}\right)$,

$$
K\left(\lambda_{d}\right) \phi=D \phi+\int_{-\tau}^{0}\left(\lambda_{d} d \mu(\theta)+d \eta(\theta)\right) e^{\lambda_{d} \theta} \int_{\theta}^{0} e^{-\lambda_{d} s} \phi(s) d s
$$

Now, we use this approach to study the asymptotic behavior of (1). Let the initial condition associated with (4) be given by

$$
X_{0}=\binom{\phi}{\phi^{\prime}} \in C\left([-\tau, 0], \mathbb{R}^{2}\right)
$$

The characteristic matrix corresponding to (4) is given by

$$
\Delta(z)=z I+z e^{-\tau z} C-E-F e^{-\tau z}=\left(\begin{array}{cc}
z & -1 \\
-q_{1}-q_{2} e^{-\tau z} & z+c z e^{-\tau z}-p_{1}-p_{2} e^{-\tau z}
\end{array}\right)
$$

so the characteristic equation is $\operatorname{det} \Delta(z)=z^{2}+c z^{2} e^{-\tau z}-\left(p_{1}+p_{2} e^{-\tau z}\right) z-q_{2} e^{-\tau z}-q_{1}$. If there exists a simple dominant zero $\lambda_{d}$ of the characteristic equation $\operatorname{det} \Delta(z)=0$, by Theorem 2.5 , we have

$$
\begin{aligned}
& \lim _{t \rightarrow \infty} e^{-\lambda_{d} t} X(t) \\
& =\lim _{t \rightarrow \infty}\binom{e^{-\lambda_{d} t} x(t)}{e^{-\lambda_{d} t} y(t)}=\left[\frac{d}{d z} \operatorname{det} \Delta\left(\lambda_{d}\right)\right]^{-1} \operatorname{adj} \Delta\left(\lambda_{d}\right) K\left(\lambda_{d}\right) \phi \\
& =\left(\begin{array}{cc}
\frac{\lambda_{d}+c \lambda_{d} e^{-\tau \lambda_{d}}-p_{1}-p_{2} e^{-\tau \lambda_{d}}}{\beta\left(\lambda_{d}\right)} & \frac{1}{\beta\left(\lambda_{d}\right)} \\
\frac{q_{1}+e_{2} e^{-\tau \lambda_{d}}}{\beta\left(\lambda_{d}\right)} & \frac{\lambda_{d}}{\beta\left(\lambda_{d}\right)}
\end{array}\right)\binom{\phi(0)}{\phi^{\prime}(0)+c \phi^{\prime}(-\tau)+\int_{-\tau}^{0}\left(p_{2}-c \lambda_{d}\right) e^{-\lambda_{d}(s+\tau)} \phi^{\prime}(s) d s+\int_{-\tau}^{0} q_{2} e^{-\lambda_{d}(s+\tau)} \phi(s) d s} .
\end{aligned}
$$

It follows that

$$
\begin{aligned}
\lim _{t \rightarrow \infty} e^{-\lambda_{d} t} x(t)= & \frac{1}{\beta\left(\lambda_{d}\right)}\left[\left(\lambda_{d}+c \lambda_{d} e^{-\tau \lambda_{d}}-p_{1}-p_{2} e^{-\tau \lambda_{d}}\right) \phi(0)+\phi^{\prime}(0)+c \phi^{\prime}(-\sigma)\right. \\
& \left.+\int_{-\tau}^{0}\left(p_{2}-c \lambda_{d}\right) e^{-\lambda_{d}(s+\tau)} \phi^{\prime}(s) d s+\int_{-\tau}^{0} q_{2} e^{-\lambda_{d}(s+\tau)} \phi(s) d s\right],
\end{aligned}
$$

where $\beta\left(\lambda_{d}\right)=2 \lambda_{d}+\left(2 c \lambda_{d}-c \tau \lambda_{d}^{2}-p_{2}+p_{2} \tau \lambda_{d}+q_{2} \tau\right) e^{-\tau \lambda_{d}}-p_{1} \neq 0$.
The next theorem gives a result similar to Theorem 2.5, in case that the real dominant eigenvalue is not simple.
Theorem 2.6. (Frasson [7]) Let $\lambda_{d}$ be a real dominant zero of $\operatorname{det} \Delta(z)$ of geometric multiplicity $n \geq 1$. If $x(t)=x(t ; \phi)$ denote the solution of (1) with initial data $x_{0}=\phi$, then the large time behaviour as a function of the initial data $\phi$ is described as follows.

1. If $P_{\lambda_{d}} \phi \neq 0$, then

$$
\lim _{t \rightarrow \infty} \frac{1}{t^{m}} e^{-\lambda_{d} t} x(t)=q_{m}\left(n, \lambda_{d}, \phi\right),
$$

where $m=\max \left\{j \in 0,1,2, \ldots n-1: q_{j}\left(n, \lambda_{d}, \phi\right) \neq 0\right\}, q_{j}$ is given by

$$
q_{j}(n, \lambda, \phi)=\frac{1}{j!} \sum_{k=j}^{n-1} \frac{D^{n-1-k} K(\lambda)}{(n-1-k)!} \frac{D_{1}^{k-j} H(\lambda, \phi)}{(k-j)!} .
$$

Furthermore, for integer $n \geq 1$, the $n-$ th Fréchet derivative of $H(\lambda, \phi)$ with respect to the first variable is given by

$$
\begin{aligned}
D_{1}^{n} H(z, \phi)= & (-1)^{n+1} n \int_{0}^{r} d \mu(\theta) \int_{0}^{\theta} \tau^{n-1} e^{-z \tau} \phi(\tau-\theta) d \tau+(-1)^{n} z \int_{0}^{r} d \mu(\theta) \int_{0}^{\theta} \tau^{n} e^{-z \tau} \phi(\tau-\theta) d \tau \\
& +(-1)^{n} \int_{0}^{r} d \eta(\theta) \int_{0}^{\theta} \tau^{n} e^{-z \tau} \phi(\tau-\theta) d \tau
\end{aligned}
$$

2. If $P_{\lambda_{d}} \phi=0$, then

$$
\lim _{t \rightarrow \infty} e^{-\lambda_{d} t} x(t)=0 .
$$

## 3. An ODE approach to asymptotic behavior

In this section, we use an ODE method to study the asymptotic behavior of the initial value problem (1). An estimate of solutions is established. As a consequence of this result, the sufficient conditions for stability, the asymptotic stability and instability of the trivial solution are presented.

The characteristic equation of (1) is

$$
\begin{equation*}
\lambda^{2}+c \lambda^{2} e^{-\lambda \tau}=p_{1} \lambda+p_{2} \lambda e^{-\lambda \tau}+q_{1}+q_{2} e^{-\lambda \tau}, \tag{8}
\end{equation*}
$$

which is obtained by seeking solutions of the form $x(t)=e^{\lambda t}$ for $t \geq-\tau$.
Suppose $\lambda_{0}$ is a real solution of the characteristic equation (8). It turns out that (1) then reduces to the first order neutral delay differential equation

$$
\begin{equation*}
z^{\prime}(t)+c e^{-\lambda_{0} \tau} z^{\prime}(t-\tau)+\left(2 \lambda_{0}-p_{1}\right) z(t)+\left(2 c \lambda_{0}-p_{2}\right) e^{-\lambda_{0} \tau} z(t-\tau)=\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \int_{-\tau}^{0} z(s+t) d s \tag{9}
\end{equation*}
$$

With (9), we associate the equation

$$
\begin{equation*}
\mu+\left(c \mu+2 c \lambda_{0}-p_{2}\right) e^{-\tau\left(\lambda_{0}+\mu\right)}+2 \lambda_{0}-p_{1}-\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \int_{-\tau}^{0} e^{\mu s} d s=0 \tag{10}
\end{equation*}
$$

which is said to be the second characteristic equation, and it is obtained from (9) by seeking solutions of the form $z(t)=e^{\mu t}$ for $t \geq-\tau$.

Now, we present a proposition, which plays a crucial role in obtaining Theorem 3.2. This proposition essentially estabishes a transformation (via a solution of the characteristic equation (8)) of the second order neutral delay differential equation (1) into the first order neutral delay differential equation (9).

Proposition 3.1. Suppose $\lambda_{0}$ is a real root of the characteristic equation (8), and let

$$
\beta\left(\lambda_{0}\right)=2 \lambda_{d}+\left(2 c \lambda_{0}-c \tau \lambda_{0}^{2}-p_{2}+p_{2} \tau \lambda_{0}+q_{2} \tau\right) e^{-\tau \lambda_{0}}-p_{1}
$$

Suppose that $\beta\left(\lambda_{0}\right) \neq 0$, then a continuous real-valued function $x$ defined on the interval $[-\tau, \infty)$ is the solution of the initial value problem (1) on $[0, \infty)$ if and only if $z$ defined by

$$
\begin{equation*}
z(t)=e^{-\lambda_{0} t} x(t)-\frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)} \quad \text { for } \quad t \geq-\tau, \tag{11}
\end{equation*}
$$

is the solution of the neutral delay differential equation (9) with the initial condition

$$
\begin{equation*}
z(t)=e^{-\lambda_{0} t} \phi(t)-\frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)} \quad \text { for } \quad-\tau \leq t \leq 0 \tag{12}
\end{equation*}
$$

where $x(t)=\phi(t)$ on $[-\tau, 0]$ and

$$
K\left(\lambda_{0}, \phi\right)=\phi^{\prime}(0)+\left(\lambda_{0}-p_{1}\right) \phi(0)+c \phi^{\prime}(-\tau)+c \lambda_{0} \phi(-\tau)-p_{2} \phi(-\tau)-\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \int_{-\tau}^{0} e^{-\lambda_{0} s} \phi(s) d s
$$

Proof. Let $x$ be the solution of the initial value problem (1) for $t \geq 0$ with $x(t)=\phi(t)$ for $-\tau \leq t \leq 0$. Define

$$
y(t)=e^{-\lambda_{0} t} x(t) \quad \text { for } \quad t \geq-\tau
$$

Using the fact that $\lambda_{0}$ is a real root of the characteristic equation (8), we have for every $t \geq 0$,

$$
\begin{align*}
{\left[y^{\prime}(t)+c e^{-\lambda_{0} \tau} y^{\prime}(t-\tau)+\left(2 \lambda_{0}-p_{1}\right) y(t)+\left(2 c \lambda_{0}-p_{2}\right) e^{-\lambda_{0} \tau} y(t-\tau)\right]^{\prime}=} & \left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) y(t) \\
& +\left(p_{2} \lambda_{0}+q_{2}-c \lambda_{0}^{2}\right) e^{-\lambda_{0} \tau} y(t-\tau) \tag{13}
\end{align*}
$$

with the initial condition satisfies

$$
\begin{equation*}
y(t)=e^{-\lambda_{0} t} \phi(t) \quad \text { for } \quad-\tau \leq t \leq 0 \tag{14}
\end{equation*}
$$

By integrating (13), and using the initial condition (14), we have

$$
\begin{equation*}
y^{\prime}(t)+c e^{-\lambda_{0} \tau} y^{\prime}(t-\tau)+\left(2 \lambda_{0}-p_{1}\right) y(t)+\left(2 c \lambda_{0}-p_{2}\right) e^{-\lambda_{0} \tau} y(t-\tau)=\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \int_{-\tau}^{0} y(s+t) d s+K\left(\lambda_{0}, \phi\right) \tag{15}
\end{equation*}
$$

for all $t \geq 0$, where $K\left(\lambda_{0}, \phi\right)$ is defined as in Proposition 3.1.
Now we suppose that $\beta\left(\lambda_{0}\right) \neq 0$ and define

$$
z(t)=y(t)-\frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)}, \quad \text { for } \quad \text { all } \quad t \geq-\tau
$$

by the definition of $\beta\left(\lambda_{0}\right)$, we obtain that $y$ satisfies (15) if and only if $z$ satisfies (11) for all $t \geq 0$. Moreover, the initial condition (14) is equivalent to (12).

An estimate of the solution of initial value problem (1) will be given in the following theorem.

Theorem 3.2. Suppose $\lambda_{0}$ is a real root of the characteristic equation (8), and let $\beta\left(\lambda_{0}\right)$ and $K\left(\lambda_{0}, \phi\right)$ be defined as in Proposition 3.1. Suppose that $\beta\left(\lambda_{0}\right) \neq 0$, let $\mu_{0}$ be a real root of the characteristic equation(10), and set

$$
\gamma\left(\lambda_{0}, \mu_{0}\right)=1+c e^{-\left(\lambda_{0}+\mu_{0}\right) \tau}-\tau\left(c \mu_{0}+2 c \lambda_{0}-p_{2}\right) e^{-\left(\lambda_{0}+\mu_{0}\right) \tau}-\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \mu_{0}^{-2}\left(\mu_{0} \tau e^{-\mu_{0} \tau}+e^{-\mu_{0} \tau}-1\right)
$$

## Define

$$
\begin{aligned}
H\left(\lambda_{0}, \mu_{0}, \phi\right)= & \phi(0)+c \phi(-\tau)+\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \int_{-\tau}^{0} e^{\mu_{0} s} \int_{s}^{0} e^{-\left(\lambda_{0}+\mu_{0}\right) u} \phi(u) d u d s \\
& -\left(c \mu_{0}+2 c \lambda_{0}-p_{2}\right) e^{-\left(\lambda_{0}+\mu_{0}\right) \tau} \int_{-\tau}^{0} e^{-\left(\lambda_{0}+\mu_{0}\right) s} \phi(s) d s \\
& -\frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)}\left[1+c e^{-\lambda_{0} \tau}+\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \mu_{0}^{-2}\left(1-e^{-\mu_{0} \tau}-\mu_{0} \tau\right)\right. \\
& \left.-\left(c \mu_{0}+2 c \lambda_{0}-p_{2}\right) \mu_{0}^{-1}\left(1-e^{-\mu_{0} \tau}\right) e^{-\left(\lambda_{0}+\mu_{0}\right) \tau}\right]
\end{aligned}
$$

We assume that the real roots $\lambda_{0}$ and $\mu_{0}$ have the following property

$$
\chi \lambda_{0}, \mu_{0}=|c| e^{-\left(\lambda_{0}+\mu_{0}\right) \tau}+\tau\left|p_{1}-\mu_{0}-2 \lambda_{0}\right|+\mu_{0}^{-2}\left(\mu_{0} \tau+e^{-\mu_{0} \tau}-1\right)\left|p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right|<1
$$

Then for any $\phi \in \mathcal{C}([-\tau, 0], \mathbb{R})$, the solution $x$ of (1) satisfies

$$
\left|e^{-\left(\mu_{0}+\lambda_{0}\right) t} x(t)-e^{-\mu_{0} t} \frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)}-\frac{H\left(\lambda_{0}, \mu_{0}, \phi\right)}{\gamma\left(\lambda_{0}, \mu_{0}\right)}\right| \leq M\left(\lambda_{0}, \mu_{0} ; \phi\right) \chi_{\lambda_{0}, \mu_{0}},
$$

where

$$
M\left(\lambda_{0}, \mu_{0} ; \phi\right)=\max _{-\tau \leq t \leq 0}\left|e^{-\mu_{0} t}\left(e^{-\lambda_{0} t} \phi(t)-\frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)}\right)-\frac{H\left(\lambda_{0}, \mu_{0} ; \phi\right)}{\gamma\left(\lambda_{0}, \mu_{0}\right)}\right|,
$$

for all $t \geq 0$.
Proof. Note that $\mu_{0}=0$ is a root of (10) if and only if $2 \lambda_{0}-p_{1}+\left(2 c \lambda_{0}-p_{2}\right) e^{-\tau \lambda_{0}}-\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \tau=0$, from the definition of $\beta\left(\lambda_{0}\right)$, we obtain that if zero is a root of (10) if and only if $\beta\left(\lambda_{0}\right)=0$. Hence, if we assume that $\beta\left(\lambda_{0}\right) \neq 0$, then we always have $\mu_{0} \neq 0$.

From the assumption that $\left|\chi_{\lambda_{0}, \mu_{0}}\right|<1$, we conclude $\gamma\left(\lambda_{0}, \mu_{0}\right)>0$. Suppose that $x$ is the solution of the initial value problem (1) with $x(\theta)=\phi(\theta)$ for $-\tau \leq \theta \leq 0$, by Proposition 3.1, the fact that $x$ is the solution of the initial value problem (1) is equivalent to the fact that $z$ is the solution of the delay differential equation (11) which satisfies the initial condition (12). Set

$$
w(t)=e^{-\mu_{0} t} z(t) \quad \text { for } \quad t \geq-\tau
$$

then by using the fact that $\mu_{0}$ is a real root of the characteristic equation (10), we obtain, for every $t \geq 0$,

$$
\begin{align*}
{\left[w(t)+c e^{-\left(\lambda_{0}+\mu_{0}\right) \tau} w(t-\tau)\right]^{\prime}=} & \left(p_{1}-\mu_{0}-2 \lambda_{0}\right) w(t)-\left(c \mu_{0}+2 c \lambda_{0}-p_{2}\right) e^{-\left(\lambda_{0}+\mu_{0}\right) \tau} w(t-\tau)  \tag{16}\\
& +\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \int_{-\tau}^{0} e^{\mu_{0} s} w(s+t) d s
\end{align*}
$$

and the initial condition satisfies

$$
\begin{equation*}
w(t)=e^{-\left(\mu_{0}+\lambda_{0}\right) t} \phi(t)-e^{-\mu_{0} t} \frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)} \quad \text { for } \quad-\tau \leq t \leq 0 . \tag{17}
\end{equation*}
$$

By integrating (16) and using the initial condition of (17), we obtain

$$
\begin{align*}
w(t)+c e^{-\left(\lambda_{0}+\mu_{0}\right) \tau} w(t-\tau)= & \left(p_{1}-\mu_{0}-2 \lambda_{0}\right) \int_{t-\tau}^{t} w(s) d s \\
& +\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \int_{-\tau}^{0} e^{\mu_{0} s} \int_{t-\tau}^{s+t} w(u) d u d s+H\left(\lambda_{0}, \mu_{0} ; \phi\right), \tag{18}
\end{align*}
$$

for all $t \geq 0$, where $H\left(\lambda_{0}, \mu_{0} ; \phi\right)$ is defined in Theorem 3.2.
Define

$$
v(t)=w(t)-\frac{H\left(\lambda_{0}, \mu_{0} ; \phi\right)}{\gamma\left(\lambda_{0}, \mu_{0}\right)} \quad \text { for } \quad \text { all } \quad t \geq-\tau
$$

then by the definition of $\gamma\left(\lambda_{0}, \mu_{0}\right)$ in Theorem 3.2, we obtain that the fact that $w$ satisfies (18) is equivalent to the fact that $v$ satisfies the following equation

$$
\begin{equation*}
v(t)+c e^{-\left(\lambda_{0}+\mu_{0}\right) \tau} v(t-\tau)=\left(p_{1}-\mu_{0}-2 \lambda_{0}\right) \int_{t-\tau}^{t} v(s) d s+\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \int_{-\tau}^{0} e^{\mu_{0} s} \int_{t-\tau}^{s+t} v(u) d u d s \tag{19}
\end{equation*}
$$

for all $t \geq 0$. Moreover, the initial condition is equivalent to

$$
\begin{equation*}
v(t)=e^{-\left(\mu_{0}+\lambda_{0}\right) t} \phi(t)-e^{-\mu_{0} t} \frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)}-\frac{H\left(\lambda_{0}, \mu_{0} ; \phi\right)}{\gamma\left(\lambda_{0}, \mu_{0}\right)}, \quad \text { for } \quad-\tau \leq t \leq 0 . \tag{20}
\end{equation*}
$$

Define

$$
M\left(\lambda_{0}, \mu_{0} ; \phi\right):=\max _{-\tau \leq \leq \leq 0}\left|e^{-\left(\mu_{0}+\lambda_{0}\right) t} \phi(t)-e^{-\mu_{0} t} \frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)}-\frac{H\left(\lambda_{0}, \mu_{0} ; \phi\right)}{\gamma\left(\lambda_{0}, \mu_{0}\right)}\right| .
$$

In view of (20), we have

$$
|v(t)| \leq M\left(\lambda_{0}, \mu_{0} ; \phi\right), \quad \text { for } \quad-\tau \leq t \leq 0 .
$$

We will next show that $M\left(\lambda_{0}, \mu_{0} ; \phi\right)$ is also a bound of $v$ on the whole positive half line. For this purpose, we take an arbitrary $\varepsilon>0$ and claim that $|v(t)|<M\left(\lambda_{0}, \mu_{0} ; \phi\right)+\varepsilon$ for $t \geq-\tau$. Indeed, suppose that there exists a point $t_{0}>0$ such that

$$
\begin{equation*}
|v(t)|<M\left(\lambda_{0}, \mu_{0} ; \phi\right)+\varepsilon, \quad \text { for } \quad-\tau \leq t<t_{0}, \quad \text { and } \quad\left|v\left(t_{0}\right)\right|=M\left(\lambda_{0}, \mu_{0} ; \phi\right)+\varepsilon . \tag{21}
\end{equation*}
$$

Then by (19) and the definition of $\chi_{\lambda_{0}, \mu_{0}}$, we have

$$
\begin{aligned}
M\left(\lambda_{0}, \mu_{0} ; \phi\right)+\varepsilon & =\left|v\left(t_{0}\right)\right| \\
& \leq|c| e^{-\left(\lambda_{0}+\mu_{0}\right) \tau}\left|v\left(t_{0}-\tau\right)\right|+\left|p_{1}-\mu_{0}-2 \lambda_{0}\right| \int_{t_{0}-\tau}^{t_{0}}|v(s)| d s+\left|p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right| \int_{-\tau}^{0} e^{\mu_{0} s} \int_{t_{0}-\tau}^{s+t}|v(u)| d u d s \\
& \leq\left(M\left(\lambda_{0}, \mu_{0} ; \phi\right)+\varepsilon\right)\left(|c| e^{-\left(\lambda_{0}+\mu_{0}\right) \tau}+\tau\left|p_{1}-\mu_{0}-2 \lambda_{0}\right|+\mu_{0}^{-2}\left(\mu_{0} \tau+e^{-\mu_{0} \tau}-1\right)\left|p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right|\right) \\
& =\left(M\left(\lambda_{0}, \mu_{0} ; \phi\right)+\varepsilon\right) \chi \lambda_{\lambda_{0}, \mu_{0}}<M\left(\lambda_{0}, \mu_{0} ; \phi\right)+\varepsilon,
\end{aligned}
$$

and we arrive at a contradiction. This implies that our claim is true and since $\varepsilon$ is arbitrary, it follows that $|v(t)| \leq$ $M\left(\lambda_{0}, \mu_{0} ; \phi\right)$ for $t \geq-\tau$. Together with (19), we arrive at

$$
\begin{aligned}
|v(t)| & \leq|c| e^{-\left(\lambda_{0}+\mu_{0}\right) \tau}|v(t-\tau)|+\left|p_{1}-\mu_{0}-2 \lambda_{0}\right| \int_{t-\tau}^{t}|v(s)| d s+\left|p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right| \int_{-\tau}^{0} e^{\mu_{0} s} \int_{t-\tau}^{s+t}|v(u)| d u d s, \\
& \leq M\left(\lambda_{0}, \mu_{0} ; \phi\right)\left(|c| e^{-\left(\lambda_{0}+\mu_{0}\right) \tau}+\tau\left|p_{1}-\mu_{0}-2 \lambda_{0}\right|+\mu_{0}^{-2}\left(\mu_{0} \tau+e^{-\mu_{0} \tau}-1\right)\left|p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right|\right) \\
& =M\left(\lambda_{0}, \mu_{0} ; \phi\right) \chi_{\lambda_{0}, \mu_{0}}<M\left(\lambda_{0}, \mu_{0} ; \phi\right),
\end{aligned}
$$

for all $t \geq 0$. This implies

$$
|v(t)|=\left|e^{-\left(\mu_{0}+\lambda_{0}\right) t} x(t)-e^{-\mu_{0} t} \frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)}-\frac{H\left(\lambda_{0}, \mu_{0} ; \phi\right)}{\gamma\left(\lambda_{0}, \mu_{0}\right)}\right| \leq M\left(\lambda_{0}, \mu_{0} ; \phi\right) \chi_{\lambda_{0}, \mu_{0}},
$$

for all $t \geq 0$. This completes the proof of Theorem 3.2.

By using the result of Theorem 3.2, we will give the asymptotic behavior of the solution of initial value problem (1) in the following Theorem.

Theorem 3.3. Suppose $\lambda_{0}$ and $\mu_{0}$ are real roots of the characteristic equations (8) and (10), respectively. Consider $\gamma\left(\lambda_{0}, \mu_{0}\right)$ and $\chi_{\lambda_{0}, \mu_{0}}$ as in Theorem 3.2. Then for any $\phi \in C([-\tau, 0], \mathbb{R})$, the solution $x$ of initial value problem (1) with $x(\theta)=\phi(\theta)$ for $-\tau \leq \theta \leq 0$ satisfies

$$
\lim _{t \rightarrow \infty} e^{-\left(\mu_{0}+\lambda_{0}\right) t} x(t)-e^{-\mu_{0} t} \frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)}=\frac{H\left(\lambda_{0}, \mu_{0} ; \phi\right)}{\gamma\left(\lambda_{0}, \mu_{0}\right)}
$$

where $K\left(\lambda_{0}, \phi\right), \beta\left(\lambda_{0}\right), H\left(\lambda_{0}, \mu_{0} ; \phi\right), \gamma\left(\lambda_{0}, \mu_{0}\right)$ are given in Proposition 3.1 and Theorem 3.2 respectively.
Proof. By the definition of $x, y, z, w$ and $v$, we have to prove that

$$
\lim _{t \rightarrow \infty} v(t)=0
$$

From Theorem 3.2, one can show by induction that $v$ satisfies

$$
\begin{equation*}
|v(t)| \leq M\left(\lambda_{0}, \mu_{0} ; \phi\right)\left(\chi_{\lambda_{0}, \mu_{0}}\right)^{n} \quad \text { for } \quad \text { all } \quad t \geq n \tau-\tau . \tag{22}
\end{equation*}
$$

Since $0 \leq \chi_{\lambda_{0}, \mu_{0}}<1$, thus from (22), it follows that $v$ tends to zero as $t \rightarrow \infty$.
Definition 3.4. The trivial solution of (1) is said to be stable iffor any $t_{0} \in \mathbb{R}$ and any $\varepsilon>0$, there exists $\delta=\delta\left(t_{0}, \varepsilon\right)>0$ such that $\left\|x_{t_{0}}\right\|<\delta$ implies $|x(t)|<\varepsilon$ for $t \geq t_{0}$. The solution is said to be asymptotically stable if it is stable and for any $t_{0} \in \mathbb{R}$ and any $\varepsilon>0$, there exists a $\delta_{a}=\delta_{a}\left(t_{0}, \varepsilon\right)>0$ such that $\left\|x_{t_{0}}\right\|<\delta_{a}$ implies $\lim _{t \rightarrow \infty} x(t)=0$.

As a consequence of Theorem 3.2 and Theorem 3.3, we have the following stability criterion.
Theorem 3.5. Let $\lambda_{0}$ and $\mu_{0}$ be real roots of the characteristic equations (8) and (10), and let $\beta\left(\lambda_{0}\right), \chi_{\lambda_{0}, \mu_{0}}, \gamma\left(\lambda_{0}, \mu_{0}\right)$ be defined as in Proposition 3.1 and Theorem 3.2 respectively, and satisfy the conditions in Theorem 3.2. Then for any $\phi \in C([-\tau, 0], \mathbb{R})$, the solution $x$ of $(1)$ with $x(\theta)=\phi(\theta)$ for $-\tau \leq \theta \leq 0$ satisfies

$$
|x(t)| \leq \frac{k_{\lambda_{0}}}{\left|\beta\left(\lambda_{0}\right)\right|} N\left(\lambda_{0}, \mu_{0} ; \phi\right) e^{\lambda_{0} t}+\left[\frac{h_{\lambda_{0}, \mu_{0}}}{\left|\gamma\left(\lambda_{0}, \mu_{0}\right)\right|}+\left(1+\frac{K_{\lambda_{0}} e_{\mu_{0}}}{\left|\beta\left(\lambda_{0}\right)\right|}+\frac{h_{\lambda_{0}, \mu_{0}}}{\left|\gamma\left(\lambda_{0}, \mu_{0}\right)\right|}\right) \chi_{\lambda_{0}, \mu_{0}}\right] N\left(\lambda_{0}, \mu_{0} ; \phi\right) e^{\left(\lambda_{0}+\mu_{0}\right) t},
$$

where

$$
\begin{aligned}
k_{\lambda_{0}}= & 1+|c|+\left|\lambda_{0}-p_{1}\right|+|c|\left|\lambda_{0}\right|+\left|p_{2}\right|+\left|p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right| \tau, \\
e_{\mu_{0}}= & \max _{-\tau \leq \leq 0}\left\{e^{-\mu_{0} t}\right\}, \\
h_{\lambda_{0}, \mu_{0}}= & 1+|c|+\left|p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right| \mu_{0}^{-2}\left(1-e^{-\mu_{0} \tau}-\mu_{0} \tau e^{-\mu_{0} \tau}\right)+\left|c \mu_{0}+2 c \lambda_{0}-p_{2}\right| \tau e^{-\left(\lambda_{0}+\mu_{0}\right) \tau} \\
& +\frac{k_{\lambda_{0}}}{\left|\beta\left(\lambda_{0}\right)\right|}\left[1+|c| e^{-\lambda_{0} \tau}+\left|p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right| \mu_{0}^{-2}\left(\mu_{0} \tau+e^{-\mu_{0} \tau}-1\right)\right. \\
& \left.+\left|c \mu_{0}+2 c \lambda_{0}-p_{2}\right|\left|\mu_{0}^{-1}\left(1-e^{-\mu_{0} \tau}\right)\right| e^{-\left(\lambda_{0}+\mu_{0}\right) \tau}\right], \\
N\left(\lambda_{0}, \mu_{0} ; \phi\right)= & \max \left\{\max _{-\tau \leq t \leq 0}\left|e^{-\lambda_{0} t} \phi(t)\right|, \max _{-\tau \leq t \leq 0}\left|e^{-\left(\lambda_{0}+\mu_{0}\right) t} \phi(t)\right|, \max _{-\tau \leq t \leq 0}\left|\phi^{\prime}(t)\right|, \max _{-\tau \leq t \leq 0}|\phi(t)|\right\} .
\end{aligned}
$$

Furthermore, the trivial solution of (1) is stable if $\lambda_{0} \leq 0, \lambda_{0}+\mu_{0} \leq 0$; it is asymptotically stable if $\lambda_{0}<0, \lambda_{0}+\mu_{0}<0$; and it is unstable if $\mu_{0}>0, \lambda_{0}+\mu_{0}>0$.
Proof. From Theorem 3.2, it follows that

$$
e^{-\left(\mu_{0}+\lambda_{0}\right) t}|x(t)| \leq \frac{\left|K\left(\lambda_{0}, \phi\right)\right|}{\left|\beta\left(\lambda_{0}\right)\right|} e^{-\mu_{0} t}+\frac{\left|H\left(\lambda_{0}, \mu_{0} ; \phi\right)\right|}{\left|\gamma\left(\lambda_{0}, \mu_{0}\right)\right|}+\left|M\left(\lambda_{0}, \mu_{0} ; \phi\right)\right| \chi_{\lambda_{0}, \mu_{0}},
$$

where $K\left(\lambda_{0}, \phi\right), H\left(\lambda_{0}, \mu_{0}, \phi\right), M\left(\lambda_{0}, \mu_{0} ; \phi\right), \beta\left(\lambda_{0}\right), \gamma\left(\lambda_{0}, \mu_{0}\right), \chi_{\lambda_{0}, \mu_{0}}$ are defined as in Theorem 3.2 respectively. From the representation of $K\left(\lambda_{0}, \phi\right), H\left(\lambda_{0}, \mu_{0}, \phi\right)$ and $M\left(\lambda_{0}, \mu_{0} ; \phi\right)$ we have

$$
\begin{aligned}
\left|K\left(\lambda_{0}, \phi\right)\right| & \leq k_{\lambda_{0}} N\left(\lambda_{0}, \mu_{0} ; \phi\right), \quad\left|H\left(\lambda_{0}, \mu_{0}, \phi\right)\right| \leq h_{\lambda_{0}, \mu_{0}} N\left(\lambda_{0}, \mu_{0} ; \phi\right), \\
\left|M\left(\lambda_{0}, \mu_{0} ; \phi\right)\right| & \leq\left(1+\frac{K_{\lambda_{0}} e_{\mu_{0}}}{\left|\beta\left(\lambda_{0}\right)\right|}+\frac{h_{\lambda_{0}, \mu_{0}}}{\left|\gamma\left(\lambda_{0}, \mu_{0}\right)\right|}\right) N\left(\lambda_{0}, \mu_{0} ; \phi\right) .
\end{aligned}
$$

Hence, it follows that

$$
e^{-\left(\mu_{0}+\lambda_{0}\right) t}|x(t)| \leq \frac{k_{\lambda_{0}}}{\left|\beta\left(\lambda_{0}\right)\right|} e^{-\mu_{0} t} N\left(\lambda_{0}, \mu_{0} ; \phi\right)+\frac{h_{\lambda_{0}, \mu_{0}}}{\left|\gamma\left(\lambda_{0}, \mu_{0}\right)\right|} N\left(\lambda_{0}, \mu_{0} ; \phi\right)+\left(1+\frac{K_{\lambda_{0}} e_{\mu_{0}}}{\left|\beta\left(\lambda_{0}\right)\right|}+\frac{h_{\lambda_{0}, \mu_{0}}}{\left|\gamma\left(\lambda_{0}, \mu_{0}\right)\right|}\right) N\left(\lambda_{0}, \mu_{0} ; \phi\right) \chi_{\lambda_{0}, \mu_{0}},
$$

which yields

$$
\begin{equation*}
|x(t)| \leq\left[\frac{h_{\lambda_{0}, \mu_{0}}}{\left|\gamma\left(\lambda_{0}, \mu_{0}\right)\right|}+\left(1+\frac{K_{\lambda_{0}} e_{\mu_{0}}}{\beta\left(\lambda_{0}\right) \mid}+\frac{h_{\lambda_{0}, \mu_{0}}}{\left|\gamma\left(\lambda_{0}, \mu_{0}\right)\right|}\right) \chi_{\lambda_{0}, \mu_{0}}\right] N\left(\lambda_{0}, \mu_{0} ; \phi\right) e^{\left(\lambda_{0}+\mu_{0}\right) t}+\frac{k_{\lambda_{0}}}{\left|\beta\left(\lambda_{0}\right)\right|} N\left(\lambda_{0}, \mu_{0} ; \phi\right) e^{\lambda_{0} t} \tag{23}
\end{equation*}
$$

for $t \geq 0$. Next, we consider three cases to discuss the stability of the trivial solution.
Case 1. We suppose that $\lambda_{0} \leq 0, \lambda_{0}+\mu_{0} \leq 0$, then $e^{\lambda_{0} t} \leq 1, e^{\left(\lambda_{0}+\mu_{0}\right) t} \leq 1$. Define $\|\phi\|=\max _{-\tau \leq t \leq 0}|\phi(t)|$, it is not difficult to obtain that $\|\phi\| \leq N\left(\lambda_{0}, \mu_{0} ; \phi\right)$. From (23), we have

$$
\begin{equation*}
|x(t)| \leq\left[\frac{k_{\lambda_{0}}}{\left|\beta\left(\lambda_{0}\right)\right|}+\left(1+\frac{K_{\lambda_{0}} e_{\mu_{0}}}{\left|\beta\left(\lambda_{0}\right)\right|}\right) \chi_{\lambda_{0}, \mu_{0}}+\left(1+\chi_{\lambda_{0}, \mu_{0}}\right) \frac{h_{\lambda_{0}, \mu_{0}}}{\left|\gamma\left(\lambda_{0}, \mu_{0}\right)\right|}\right] N\left(\lambda_{0}, \mu_{0} ; \phi\right) \tag{24}
\end{equation*}
$$

for every $t \geq 0$. Define

$$
\rho:=\frac{k_{\lambda_{0}}}{\left|\beta\left(\lambda_{0}\right)\right|}+\left(1+\frac{K_{\lambda_{0}} e_{\mu_{0}}}{\left|\beta\left(\lambda_{0}\right)\right|}\right) \chi_{\lambda_{0}, \mu_{0}}+\left(1+\chi_{\lambda_{0}, \mu_{0}}\right) \frac{h_{\lambda_{0}, \mu_{0}}}{\left|\gamma\left(\lambda_{0}, \mu_{0}\right)\right|} .
$$

For any $\varepsilon>0$, we choose $\delta=\varepsilon \rho^{-1}$ such that $N\left(\lambda_{0}, \mu_{0} ; \phi\right)<\delta$, since $\|\phi\| \leq N\left(\lambda_{0}, \mu_{0} ; \phi\right)$, we obtain that $\|\phi\| \leq \delta$. From estimate (24), we obtain $|x(t)| \leq \rho N\left(\lambda_{0}, \mu_{0} ; \phi\right)<\rho \delta=\varepsilon$. This implies the trivial solution of (1) is stable.

Case 2. We suppose that $\lambda_{0}<0, \lambda_{0}+\mu_{0}<0$. From estimate (23), it follows that $\lim _{t \rightarrow \infty} x(t)=0$. Hence, the trivial solution of (1) is asymptotically stable.

Case 3. Let $\mu_{0}>0, \lambda_{0}+\mu_{0}>0$. If the trivial solution of (1) is stable, then there exists a number $l=l(1)>0$ such that, for any $\phi \in C([-\tau, 0], \mathbb{R})$ with $\|\phi\|<l$ the solution $x$ of (1) with $x(\theta)=\phi(\theta)$ for $-\tau \leq \theta \leq 0$ satisfies $|x(t)|<1$ for $t \geq 0$. Define

$$
\phi_{0}(t)=e^{\left(\lambda_{0}+\mu_{0}\right) t}-e^{\lambda_{0} t} \quad \text { for } t \in[-\tau, 0]
$$

By definition of $K\left(\lambda_{0}, \phi\right)$ and $H\left(\lambda_{0}, \mu_{0}, \phi\right)$, and using the relation of (8), we have that $K\left(\lambda_{0}, \phi_{0}\right)=-\beta\left(\lambda_{0}\right)$ and $H\left(\lambda_{0}, \mu_{0}, \phi_{0}\right)=\gamma\left(\lambda_{0}, \mu_{0}\right)$. Let $\phi \in C([-\tau, 0], \mathbb{R})$ be defined by $\phi=\frac{l_{1}}{\left\|\phi_{0}\right\|} \phi_{0}$ with $0<l_{1}<l$. From Theorem 3.3, we have

$$
\begin{equation*}
\lim _{t \rightarrow \infty} e^{-\left(\mu_{0}+\lambda_{0}\right) t} x(t)-e^{-\mu_{0} t} \frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)}=\frac{H\left(\lambda_{0}, \mu_{0} ; \phi\right)}{\gamma\left(\lambda_{0}, \mu_{0}\right)} . \tag{25}
\end{equation*}
$$

On the other hand,

$$
\lim _{t \rightarrow \infty} e^{-\left(\mu_{0}+\lambda_{0}\right) t} x(t)-e^{-\mu_{0} t} \frac{K\left(\lambda_{0}, \phi\right)}{\beta\left(\lambda_{0}\right)}=\lim _{t \rightarrow \infty} e^{-\left(\mu_{0}+\lambda_{0}\right) t} x(t)+\frac{l_{1}}{\left\|\phi_{0}\right\|} e^{-\mu_{0} t}=0
$$

but

$$
\frac{H\left(\lambda_{0}, \mu_{0} ; \phi\right)}{\gamma\left(\lambda_{0}, \mu_{0}\right)}=\frac{\left(l_{1} /\left\|\phi_{0}\right\|\right) H\left(\lambda_{0}, \mu_{0} ; \phi\right)}{\gamma\left(\lambda_{0}, \mu_{0}\right)}=\frac{l_{1}}{\left\|\phi_{0}\right\|}>0 .
$$

This is a contradiction to (25) and this shows that the trivial solution of (1) is unstable.

## 4. Discussion of the two approaches

For the initial value problem (3), we first consider the scalar case $C=C([-\tau, 0], \mathbb{R})$. The characteristic equation $\Delta(z)$ is given by (6). Define the auxiliary function $\chi: \mathbb{C} \rightarrow[0, \infty)$ by

$$
\begin{equation*}
\chi(z)=\int_{-\tau}^{0}(1-\theta|z|)\left|e^{z \theta}\right| d V(\mu)(\theta)+\int_{-\tau}^{0}(-\theta)\left|e^{z \theta}\right| d V(\eta)(\theta) \tag{26}
\end{equation*}
$$

where $V(\mu)(\theta)$ denotes the total variation function of $\mu$ on $[-\tau, \theta]$ for each $\theta$ in $(-\tau, 0]$.
Theorem 4.1. (Frasson [6]) Suppose that $z_{0} \in \mathbb{C}$ is a zero of $\operatorname{det} \Delta(z)$ in (6). If $\chi\left(z_{0}\right)<1$, then $z_{0}$ is a simple dominant zero of $\Delta(z)$.

Combining this result with Theorem 2.5, we arrive at
Theorem 4.2. Let $x(\cdot)$ be the solution of (3) subjected to the initial condition $x_{0}=\phi \in \mathcal{C}([-\tau, 0], \mathbb{R})$. If $\lambda_{d}$ is a real zero of characteristic equation $\Delta(z)$ given by (6) such that $\chi\left(\lambda_{d}\right)<1$, where $\chi(\cdot)$ is given by (26), then the asymptotic behavior of $x(\cdot)$ is given by

$$
\lim _{t \rightarrow \infty} e^{-\lambda_{d} t} x(t)=\frac{1}{H\left(\lambda_{d}\right)} K\left(\lambda_{d}\right) \phi,
$$

where

$$
\begin{aligned}
H\left(\lambda_{d}\right) & =1-\int_{-\tau}^{0} e^{\lambda_{d} \theta} d \mu(\theta)-\int_{-\tau}^{0} \theta e^{\lambda_{d} \theta}\left(\lambda_{d} d \mu(\theta)+d \eta(\theta)\right) \\
K\left(\lambda_{d}\right) \phi & =M \psi+\int_{-\tau}^{0}\left(\lambda_{d} d \mu(\theta)+d \eta(\theta)\right) e^{\lambda_{d} \theta} \int_{\theta}^{0} e^{-\lambda_{d} s} \psi(s) d s
\end{aligned}
$$

On the other side, by using the ODE approach (see [15]), we arrive at the same conclusion as in Theorem 4.2. This means that in this case the spectral approach is equivalent to the ODE appraoch.

## Example 4.3.

$$
\left\{\begin{array}{l}
x^{\prime}(t)+c x^{\prime}(t-\sigma)=a x(t)+b x(t-\tau)  \tag{27}\\
x(t)=\phi(t), \quad-\tau \leq t \leq 0
\end{array}\right.
$$

The characteristic equation of (27) is

$$
\begin{equation*}
\Delta(\lambda)=\lambda\left(1+c e^{-\lambda \sigma}\right)-a-b e^{-\lambda \tau} \tag{28}
\end{equation*}
$$

By applying the ODE approach, Kordonis, Niyianni and Philos [12] obtained the following theorem.
Theorem 4.4. ([12]) Let $\lambda_{0}$ be a real root of characteristic equation (28) with the property $|c|\left(1+\left|\lambda_{0}\right| \sigma\right) e^{-\lambda_{0} \sigma}+|b| \tau e^{-\lambda_{0} \tau}<1$, then for any $\phi \in \mathcal{C}([-\tau, 0], \mathbb{R})$, we have

$$
\lim _{t \rightarrow \infty} e^{-\lambda_{0} t} x(\phi ; t)=\frac{L\left(\lambda_{0} ; \phi\right)}{1+\gamma\left(\lambda_{0}\right)}
$$

where

$$
\begin{aligned}
\gamma\left(\lambda_{0}\right) & =c\left(1-\lambda_{0} \sigma\right) e^{-\lambda_{0} \sigma}+b \tau e^{-\lambda_{0} \tau} \\
L\left(\lambda_{0} ; \phi\right) & =\phi(0)+c \phi(-\sigma)-c \lambda_{0} e^{-\lambda_{0} \sigma} \int_{-\sigma}^{0} e^{\lambda_{0} s} \phi(s) d s+b e^{-\lambda_{0} \tau} \int_{-\tau}^{0} e^{\lambda_{0} s} \phi(s) d s
\end{aligned}
$$

On the other hand, if $\lambda_{0}$ is a real root of characteristic equation (28) and has the property as in Theorem 4.4, by Theorem 4.1, $\lambda_{0}$ is a simple dominant root of (28). Hence, applying Theorem 4.2, we have

$$
\lim _{t \rightarrow \infty} e^{-\lambda_{0} t} x(\phi ; t)=\frac{K\left(\lambda_{0} ; \phi\right)}{H\left(\lambda_{0}\right)},
$$

where

$$
\begin{aligned}
K\left(\lambda_{0} ; \phi\right) & =\phi(0)+c \phi(-\sigma)-c \lambda_{0} e^{-\lambda_{0} \sigma} \int_{-\sigma}^{0} e^{\lambda_{0} s} \phi(s) d s+b e^{-\lambda_{0} \tau} \int_{-\tau}^{0} e^{\lambda_{0} s} \phi(s) d s \\
H\left(\lambda_{0}\right) & =1+c\left(1-\lambda_{0} \sigma\right) e^{-\lambda_{0} \sigma}+b \tau e^{-\lambda_{0} \tau}
\end{aligned}
$$

which is the same as the result in Theorem 4.4.
Next, we consider the conditions of Theorem 3.3 in more detail. Suppose $\mu_{0}$ is a real root of the second characteristic equation (10). If $\mu_{0}$ satisfies $\chi_{\lambda_{0}, \mu_{0}}<1$, we claim that $\mu_{0}$ is a simple dominant zero. Let

$$
G(\mu):=\mu+\left(c \mu+2 c \lambda_{0}-p_{2}\right) e^{-\tau\left(\lambda_{0}+\mu\right)}+2 \lambda_{0}-p_{1}-\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \int_{-\tau}^{0} e^{\mu s} d s
$$

By the condition $\chi_{\lambda_{0}, \mu_{0}}<1$ in Theorem 3.3, we have $G^{\prime}\left(\mu_{0}\right) \neq 0$.
Indeed, since $\chi_{\lambda_{0}, \mu_{0}}<1$,

$$
\left|G^{\prime}\left(\mu_{0}\right)\right| \geq 1-\left[|c| e^{-\left(\lambda_{0}+\mu_{0}\right) \tau}+\tau\left|p_{1}-\mu_{0}-2 \lambda_{0}\right|+\mu_{0}^{-2}\left(\mu_{0} \tau+e^{-\mu_{0} \tau}-1\right)\left|p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right|\right]>0 .
$$

Since $\chi_{\lambda_{0}, \mu_{0}}<1$, let $0<\delta<1$ such that $\chi_{\lambda_{0}, \mu_{0}}<\delta$. From the representation of $\chi_{\lambda_{0}, \mu_{0}}$, we can estimate

$$
1-\frac{1}{\delta}|c| e^{-\left(\lambda_{0}+\mu_{0}\right) \tau}>\frac{1}{\delta}\left\{\tau\left|p_{1}-\mu_{0}-2 \lambda_{0}\right|+\mu_{0}^{-2}\left(\mu_{0} \tau+e^{-\mu_{0} \tau}-1\right)\left|p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right|\right\}
$$

Let $\varepsilon>0$ such that $1<e^{\varepsilon r} \leq \frac{1}{\delta}$, and we let $\Omega$ denote the right half plane given by

$$
\Omega=\left\{\mu \in \mathbb{C}: \operatorname{Re} \mu>\mu_{0}-\varepsilon\right\}
$$

For $\mu \in \Omega$ and $0 \leq s \leq r$, we have $\left|e^{-s \mu}\right|=e^{-s R e \mu}<e^{-s \mu_{0}} e^{\varepsilon s} \leq \frac{e^{-s \mu_{0}}}{\delta}$. If $\mu \in \Omega$, let $\gamma$ denote the line segment between $\mu_{0}$ and $\mu$, such that the segment is in $\Omega$, then for $0 \leq s \leq r$,

$$
\begin{equation*}
\left|e^{-s \mu}-e^{-s \mu_{0}}\right|=\left|\int_{\mu_{0}}^{\mu} s e^{-t s} d t\right|=s\left|\int_{\gamma} e^{-t s} d t\right| \leq \frac{e^{-s \mu_{0}}}{\delta}\left|\mu-\mu_{0}\right| s, \tag{29}
\end{equation*}
$$

since $G\left(\mu_{0}\right)=0$, we have

$$
\begin{aligned}
G(\mu)= & \left(\mu-\mu_{0}\right)\left(1+c e^{-\tau\left(\lambda_{0}+\mu\right)}\right)+c \mu_{0} e^{-\tau \lambda_{0}}\left(e^{-\tau \mu}-e^{-\tau \mu_{0}}\right) \\
& +\left(2 c \lambda_{0}-p_{2}\right) e^{-\tau \lambda_{0}}\left(e^{-\tau \mu}-e^{-\tau \mu_{0}}\right)-\left(p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right) \int_{-\tau}^{0}\left(e^{-\tau \mu}-e^{-\tau \mu_{0}}\right) d s .
\end{aligned}
$$

Now, we estimate $|G(\mu)|$ by using (29),

$$
|G(\mu)| \geq\left|\mu-\mu_{0}\right|\left(1-\left|\frac{1}{\delta} c e^{-\tau\left(\lambda_{0}+\mu_{0}\right)}\right|\right)-\frac{\left|\mu-\mu_{0}\right|}{\delta}\left\{\tau\left|p_{1}-\mu_{0}-2 \lambda_{0}\right|+\mu_{0}^{-2}\left(\mu_{0} \tau+e^{-\mu_{0} \tau}-1\right)\left|p_{1} \lambda_{0}+q_{1}-\lambda_{0}^{2}\right|\right\}>0
$$

which means $\mu_{0}$ is the only zero in the right half plane $\Omega$, so $\mu_{0}$ is a simple dominant zero.
For the case when the space $C=C\left([-\tau, 0], \mathbb{R}^{n}\right)$, as we discussed in Section 2, the spectral approach can be applied to study the asymptotic behavior of the functional differential equations with the solutions in this space. However, for this case, the ODE approach is not applicable.

In the following, we present two examples to illustrate the relations of the two approaches.

Example 4.5. If $a=1, b=1, c=1$ and $\sigma=\tau=1$ in (1), then we have

$$
\left\{\begin{array}{l}
x^{\prime \prime}(t)+x^{\prime \prime}(t-1)=x(t)+x(t-1)  \tag{30}\\
x(t)=\phi(t), \quad-1 \leq t \leq 0
\end{array}\right.
$$

The characteristic equation of (30) is $\lambda^{2}+\lambda^{2} e^{-\lambda}=1+e^{-\lambda}$. We denote $F_{1}(\lambda)=\lambda^{2}+\lambda^{2} e^{-\lambda}-1-e^{-\lambda}=\left(\lambda^{2}-1\right)\left(1+e^{-\lambda}\right)$. Since $F_{1}(1)=0, F_{1}^{\prime}(1)=2+\frac{2}{e} \neq 0$, we have that $\lambda_{0}=1$ is a simple zero of $F_{1}(\lambda)$. Hence, (30) becomes

$$
\begin{equation*}
z^{\prime}(t)+e^{-1} z^{\prime}(t-1)+2 z(t)+2 e^{-1} z(t-1)=0 \tag{31}
\end{equation*}
$$

and the characteristic equation of (31) is

$$
\mu+(\mu+2) e^{-(\mu+1)}+2=(\mu+2)\left(1+e^{-\mu-1}\right)=0
$$

We denote $G_{1}(\mu)=(\mu+2)\left(1+e^{-\mu-1}\right)$. Since $\mu_{0}=-2$ is a real zero of $G_{1}(\mu)$, the condition of Theorem 3.3 is $\chi_{1,-2}=e>1$, so Theorem 3.3 is not applicable.

But $\lambda=-1$ is another root of $F_{1}(\lambda)$ and satisfies $F_{1}^{\prime}(-1)=-2-2 e \neq 0$, so (30) becomes

$$
\begin{equation*}
z^{\prime}(t)+e z^{\prime}(t-1)-2 z(t)-2 e z(t-1)=0 \tag{32}
\end{equation*}
$$

and the characteristic equation of (32) is

$$
\begin{equation*}
\mu+(\mu-2) e^{-(\mu-1)}-2=(\mu-2)\left(1+e^{-(\mu-1)}\right)=0 \tag{33}
\end{equation*}
$$

It is easy to check that $\mu=\mu_{0}=2$ is a real root of (33). Corresponding to the roots $\lambda_{0}=-1$ and $\mu_{0}=2$, the condition of Theorem 3.3 becomes $\chi_{-1,2}=e^{-1}<1$. Therefore by using the result of Theorem 3.3, the asymptotic behavior of initial value problem (30) is

$$
\lim _{t \rightarrow \infty} e^{-t} x(t)=\frac{H(-1,2 ; \phi)}{\gamma(-1,2)}=\frac{\phi(0)+\phi(-1)+\phi^{\prime}(0)+\phi^{\prime}(-1)}{2+2 e^{-1}}
$$

Next, we apply Theorem 2.5 to study the asymptotic behavior of initial value problem (30). The characteristic matrix of (30) is

$$
\Delta(z)=\left(\begin{array}{cc}
z & -1 \\
-e^{-z}-1 & z+z e^{-z}
\end{array}\right)
$$

Since $z=z_{0}=1$ is a dominant zero of det $\Delta(z)$, and $\left.\frac{d}{d z}(\operatorname{det} \Delta(z))\right|_{z=z_{0}}=2+2 e^{-1} \neq 0$, we obtain that $z_{0}=1$ is a simple dominant zero of $\operatorname{det} \Delta(z)$, which satisfies the condition of Theorem 2.5. Therefore, we have

$$
\lim _{t \rightarrow \infty} e^{-t} x(t)=\frac{\phi(0)+\phi(-1)+\phi^{\prime}(0)+\phi^{\prime}(-1)}{2+2 e^{-1}}
$$

From this example, we see that the result by the spectral approach is the same as the one by the ODE approach.
Example 4.6. We suppose $a=1, \sigma=\tau=1, b=c$ in (1), we have

$$
\left\{\begin{array}{l}
x^{\prime \prime}(t)+c x^{\prime \prime}(t-1)=x(t)+c x(t-1)  \tag{34}\\
x(t)=\phi(t), \quad-1 \leq t \leq 0
\end{array}\right.
$$

The characteristic equation of (34) is

$$
\lambda^{2}+c \lambda^{2} e^{-\lambda}=1+c e^{-\lambda}
$$

we denotes $F_{2}(\lambda)=\lambda^{2}+c \lambda^{2} e^{-\lambda}-1-c e^{-\lambda}=\left(\lambda^{2}-1\right)\left(1+c e^{-\lambda}\right)$. Since $F_{2}(-1)=0, F_{2}^{\prime}(-1)=-2-2 c e \neq 0$, So $\lambda_{0}=-1$ is a simple zero of $F_{2}(\lambda)$, (34) becomes

$$
\begin{equation*}
z^{\prime}(t)+c e z^{\prime}(t-1)-2 z(t)-2 c e z(t-1)=0 \tag{35}
\end{equation*}
$$

The characteristic equation of (35) is

$$
\mu+(\mu-2) c e^{-(\mu-1)}-2=(\mu-2)\left(1+c e^{-(\mu-1)}\right)=0 .
$$

We denote $G_{2}(\mu)=(\mu-2)\left(1+c e^{-(\mu-1)}\right)$. Since $\mu_{0}=2$ is a real zero of $G_{2}(\mu)$, corresponding to the roots $\lambda_{0}=-1$ and $\mu_{0}=2$, the condition of Theorem 3.3 is $\chi_{-1,2}=|c| e^{-1}$. If $|c|<e$, we have $\chi_{-1,2}<1$. Therefore by using the result of Theorem 3.3, the asymptotic behavior of initial value problem (34) is

$$
\lim _{t \rightarrow \infty} e^{-t} x(t)=\frac{H(-1,2 ; \phi)}{\gamma(-1,2)}=\frac{\phi(0)+\phi^{\prime}(0)+c\left(\phi(-1)+\phi^{\prime}(-1)\right)}{2+2 c e^{-1}} .
$$

Next, we consider (34) by applying spectral approach. For (34), the characteristic matrix is given by

$$
\Delta(z)=\left(\begin{array}{cc}
z & -1 \\
-c e^{-z}-1 & z+c z e^{-z}
\end{array}\right)
$$

(1) Case $-e<c$. It is not difficult to check $z_{0}=1$ is a dominant zero of $\operatorname{det} \Delta(z)$, since $\left.\frac{d}{d z} \operatorname{det} \Delta(z)\right|_{z=z_{0}}=2+2 c e^{-1} \neq$ 0 , so $z_{0}=1$ is a simple dominant zero of $\operatorname{det} \Delta(z)$. Therefore, by applying the result of Theorem 2.5,

$$
\lim _{t \rightarrow \infty} e^{-t} x(t)=\frac{\phi(0)+\phi^{\prime}(0)+c\left(\phi(-1)+\phi^{\prime}(-1)\right)}{2+2 c e^{-1}} .
$$

(2) Case $c<-e$. After checking the roots of $\operatorname{det} \Delta(z)$, we find $z_{0}=\ln (-c)$ is a dominant zero of det $\Delta(z)$ and we can also use Theorem 2.5 to obtain the asymptotic behavior of the equation initial value problem (34).
(3) Case $c=-e$. We learned that $z_{0}=1$ is a dominant zero with order 2 , so by the spectral approach in [6], we can have the asymptotic behavior of the equation initial value problem (34).
From this example, we derive that for the ODE approach, the coefficient $c$ should satisfy $|c|<e$. However, for every $c \in \mathbb{R}$, the asymptotic behavior of the equation initial value problem (34) can be obtained by the spectral projection approach, and the result is the same as the one by the ODE approach when $c$ satisfies $|c|<e$.

## 5. A fixed point approach towards asymptotic behavior

In this section, we study the special case of the system (1) with $c=0$ and $p_{2}=0$. Since it is not easy to apply the ODE approach or the spectral approach to discuss its asymptotic behavior, we introduce a third approach. This approach is based on fixed point theory and relies on three principles: a complete metric space, the contraction mapping principle, and an elementary variation of parameters formula. Together this yields existence, uniqueness and stability.

By using a fixed point approach, Burton and Furumochi [1] have considered asymptotic behavior of solutions of the following linear equation

$$
\begin{equation*}
x^{\prime \prime}(t)+a x^{\prime}(t)+b x(t-r)=0 \tag{36}
\end{equation*}
$$

and obtained the following.
Theorem 5.1. (Burton and Furumochi [1]) Let $a>0$ and $b>0$. If

$$
\operatorname{br}\left(1+\int_{0}^{t}\left|A e^{A(t-s)}\right| d s\right)<1
$$

holds, where $A=\left(\begin{array}{cc}0 & 1 \\ -b & -a\end{array}\right)$, then every solution of equation (36) and its derivative tend to 0 as $t \rightarrow \infty$.

By using a similar technique as introduced by Burton and Furumochi [1], we consider the retarded delay differential equation

$$
\begin{equation*}
x^{\prime \prime}(t)+a x^{\prime}(t)+b x(t-r)+c x(t)=0 \tag{37}
\end{equation*}
$$

If we set $x^{\prime}=y$, then (37) can be written in the following form

$$
y^{\prime}=-a y-(b+c) x+(d / d t) \int_{t-r}^{t} b x(s) d s
$$

which is then expressed as the vector system

$$
z^{\prime}=A z+(d / d t) \int_{t-r}^{t} B z(s) d s
$$

where $A$ and $B$ are

$$
A=\left(\begin{array}{cc}
0 & 1  \tag{38}\\
-(b+c) & -a
\end{array}\right) \quad \text { and } \quad B=\left(\begin{array}{ll}
0 & 0 \\
b & 0
\end{array}\right)
$$

By the variation of parameters formula, we arrive at the following representation

$$
z(t)=e^{A t} z_{0}+\int_{0}^{t} e^{A(t-s)}(d / d s) \int_{s-r}^{s} B z(u) d u d s
$$

employing an integration by parts, we have

$$
z(t)=e^{A t} z_{0}+\int_{t-r}^{t} B z(u) d u-e^{A t} \int_{-r}^{0} B z(u) d u+A \int_{0}^{t} e^{A(t-s)}(d / d s) \int_{s-r}^{s} B z(u) d u d s
$$

In order to have $e^{A t} \rightarrow 0$ as $t \rightarrow \infty$, we need

$$
\begin{equation*}
b+c>0, \quad a>0 \tag{39}
\end{equation*}
$$

Let $C=C\left([-r, 0], \mathbb{R}^{2}\right)$ be the space of continuous functions and $\phi \in C$ be a given initial function. Set

$$
S_{\phi}:=\left\{\varphi: \varphi \in C\left([-r, \infty), \mathbb{R}^{2}\right), \varphi(t)=\phi(t) \quad \text { on }[-r, 0], \quad \varphi(t) \rightarrow 0 \quad \text { as } t \rightarrow \infty\right\} .
$$

For $\varphi \in S_{\phi}$, define the operator $P$ by

$$
(P \varphi)(t)=e^{A t} \phi(0)+\int_{t-r}^{t} B \varphi(u) d u-e^{A t} \int_{-r}^{0} B \varphi(u) d u+A \int_{0}^{t} e^{A(t-s)}(d / d s) \int_{s-r}^{s} B \varphi(u) d u d s
$$

Next we introduce a norm. If

$$
z=\binom{x}{y}
$$

then define $|z|_{0}:=|x|+|y|$ to be the $l_{1}$-norm. Let $Q$ be a fixed $2 \times 2$ nonsingular matrix such that $|q|_{0} \leq 1$, where $q$ denotes the second column of $Q$, and define $|z|:=|Q z|_{0}$. For a $2 \times 2$ matrix $M$, define

$$
|M|:=\sup \left\{\left|Q M Q^{-1} z\right|_{0}:|z|_{0}=1\right\}
$$

then $|M|$ is the norm of $M$. Next we can formulate the following theorem.
Theorem 5.2. Let $b+c>0, b>0$ and $a>0$. If the following condition is satisfied

$$
\begin{equation*}
b r\left(1+\int_{0}^{t}\left|A e^{A(t-s)}\right| d s\right)<1 \tag{40}
\end{equation*}
$$

where $A$ is given by (38), then every solution of equation (37) and its derivative tend to 0 as $t \rightarrow \infty$.

Proof. Since $e^{A t}$ is a $L^{1}$-function on $\mathbb{R}^{+}$, if $\varphi \in S_{\phi}$, then $(P \varphi)(t) \rightarrow 0$ as $t \rightarrow \infty$. Thus, $P: S_{\phi} \rightarrow S_{\phi}$. Furthermore, from condition (40), we have that $P$ is a contraction. The proof proceeds similarly to the proof as presented by Burton and Furumochi [1].

Example 5.3. Consider the equation

$$
\begin{equation*}
x^{\prime \prime}+\frac{7}{12} x^{\prime}+\frac{1}{6} x(t-2)-\frac{1}{12} x=0, \quad t \in \mathbb{R}^{+} \tag{41}
\end{equation*}
$$

We have

$$
A=\left(\begin{array}{cc}
0 & 1 \\
-1 / 12 & -7 / 12
\end{array}\right) \quad \text { and } \quad B=\left(\begin{array}{cc}
0 & 0 \\
1 / 6 & 0
\end{array}\right)
$$

The eigenvalues of $A$ are $-\frac{1}{3}$ and $-\frac{1}{4}$, let $Q$ be a $2 \times 2$ nonsingular matrix such that

$$
Q A Q^{-1}=\left(\begin{array}{cc}
-1 / 3 & 0 \\
0 & -1 / 4
\end{array}\right)
$$

Then we have $A e^{A(t-s)}=Q E Q^{-1}$, where

$$
E=\left(\begin{array}{cc}
-(1 / 3) e^{-(t-s) / 3} & 0 \\
0 & -(1 / 4) e^{-(t-s) / 4}
\end{array}\right)
$$

and

$$
\begin{aligned}
\left|A e^{A(t-s)}\right| & =\sup \left\{|E z|_{0}:|z|_{0}=1\right\} \\
& =\sup \left\{(|x| / 3) e^{-(t-s) / 3}+(|y| / 4) e^{-(t-s) / 4}:|x|+|y|=1\right\} \\
& \leq(1 / 3) e^{-(t-s) / 3}+(1 / 4) e^{-(t-s) / 4}
\end{aligned}
$$

Hence,

$$
\int_{0}^{t}\left|A e^{A(t-s)}\right| d s \leq \int_{0}^{t}\left[(1 / 3) e^{-(t-s) / 3}+(1 / 4) e^{-(t-s) / 4}\right] d s=2-e^{-t / 3}-e^{-t / 4}<2, \quad t \geq 0
$$

which together with $b r=\frac{1}{3}$ implies that (40) holds. Thus, by Theorem 5.2, we have that every solution of (41) and its derivative tends to 0 as $t \rightarrow \infty$. An extensive discussion of the fixed point approach to determine stability and asymptotic behaviour of various types of (neutral) delay differential equations can be found in [2].

## 6. Application to a mechanical model of turning processes

Systems governed by (neutral) delay differential equations (DDEs) often come up in different fields of science and engineering. An important mechanical application are turning processes. For the simplest model of turning, the governing equation of motion is an autonomous DDE with a corresponding infinite dimensional state space, see $[10,11]$. This fact results in an infinite number of complex characteristic roots, most of them having negative real parts referring to damped components of the vibration signals. There are finitely many characteristic roots with positive real part.

From the detailed introduction of mechanical models of turning processes in [10], we focus on a linear autonomous delay differential equation

$$
\begin{equation*}
m \xi^{\prime \prime}(t)+c \xi^{\prime}(t)+k \xi(t)=-w h(\xi(t)-\xi(t-\tau)) \tag{42}
\end{equation*}
$$

where $m, c, k, w, h, \tau$ are constants. For the meaning of every parameter, refer to [10].
Using the model parameters, equation (42) reads

$$
\begin{equation*}
\xi^{\prime \prime}(t)+2 \zeta w_{n} \xi^{\prime}(t)+w_{n}^{2} \xi(t)=-\frac{w h}{m}(\xi(t)-\xi(t-\tau)) \tag{43}
\end{equation*}
$$



Figure 1: The solution of equation (42).
where $w_{n}=\sqrt{k / m}, \zeta=c /\left(2 m w_{n}\right)$. Generally, $\zeta \approx 0.005 \sim 0.02$. Equation (43) is the standard linear delay differential equation model of the turning process.

Equation (43) can be even further simplified. Introduce the dimensionless time $\tilde{t}$ by $\tilde{t}=t w_{n}$, and by abuse of notation, drop the tilde immediately. This gives the dimensionless equation of motion

$$
\begin{equation*}
\xi^{\prime \prime}(t)+2 \zeta \xi^{\prime}(t)+\xi(t)=-\tilde{w}\left(\xi(t)-\xi\left(t-w_{n} \tau\right)\right) \tag{44}
\end{equation*}
$$

where $\tilde{w}=\frac{w h}{m w_{n}^{2}}$. In the following, we study the asymptotic stability of equation (44) by Theorem 5.2.
Equation (44) can be written as the following.

$$
\begin{equation*}
\xi^{\prime \prime}(t)+2 \zeta \xi^{\prime}(t)+(1+\tilde{w}) \xi(t)-\tilde{w} \xi\left(t-w_{n} \tau\right)=0 \tag{45}
\end{equation*}
$$

We denote $w_{n} \tau=r$. From (45), we have

$$
A=\left(\begin{array}{cc}
0 & 1 \\
-1 & -2 \zeta
\end{array}\right) \quad \text { and } \quad B=\left(\begin{array}{cc}
0 & 0 \\
-\tilde{w} & 0
\end{array}\right)
$$

The characteristic equation of $A$ is

$$
\lambda^{2}+2 \zeta \lambda+1=0
$$

The eigenvalues are

$$
\lambda_{1}=-\zeta+i \sqrt{1-\zeta^{2}}, \quad \lambda_{2}=-\zeta-i \sqrt{1-\zeta^{2}}
$$

so $\left|\lambda_{1}\right|=\left|\lambda_{2}\right|=1,\left|\lambda_{1} e^{\lambda_{1}(t-s)}\right|=\left|\lambda_{1}\right| e^{\left(R e \lambda_{1}\right)(t-s)}$ and $\left|\lambda_{2} e^{\lambda_{2}(t-s)}\right|=\left|\lambda_{2}\right| e^{\left(\text {Re } \lambda_{2}\right)(t-s)}$.
If $\zeta \notin\{-1,1\}$, the two different eigenvalues $\lambda_{1}, \lambda_{2}$ have eigenvectors $V_{1}$ and $V_{2}$, which are linearly independent. We suppose that $Q=\left(V_{1}, V_{2}\right)^{-1}$, then

$$
Q A Q^{-1}=\left(\begin{array}{cc}
\lambda_{1} & 0 \\
0 & \lambda_{2}
\end{array}\right) \triangleq \Lambda
$$

Hence,

$$
A=Q^{-1} \Lambda Q=\left(V_{1}, V_{2}\right) A\left(V_{1}, V_{2}\right)^{-1}, \quad e^{A(t-s)}=Q^{-1} e^{\Lambda(t-s)} Q
$$

Then we have

$$
A e^{A(t-s)}=Q^{-1} \Lambda Q Q^{-1} e^{\Lambda(t-s)} Q=Q^{-1} \Lambda e^{\Lambda(t-s)} Q=Q^{-1} E Q,
$$

where

$$
E=\left(\begin{array}{cc}
\lambda_{1} e^{\lambda_{1}(t-s)} & 0 \\
0 & \lambda_{2} e^{\lambda_{2}(t-s)}
\end{array}\right) .
$$

Using the norm in Theorem 5.2, we have

$$
\begin{aligned}
\left|A e^{A(t-s)}\right| & =\sup \left\{|E z|_{0}:|z|_{0}=1\right\} \\
& =\sup \left\{\left|x \lambda_{1} e^{\lambda_{1}(t-s)}\right|+\left|y \lambda_{2} e^{\lambda_{2}(t-s)}\right|:|x|+|y|=1\right\} \\
& =\sup \left\{|x|\left|\lambda_{1}\right| e^{\left(R e \lambda_{1}\right)(t-s)}+|y|\left|\lambda_{2}\right| e^{\left(\text {Red } \lambda_{2}\right)(t-s)}:|x|+|y|=1\right\} \\
& \leq 2 e^{-\zeta(t-s)} .
\end{aligned}
$$

Hence,

$$
\begin{gathered}
\int_{0}^{t}\left|A e^{A(t-s)}\right| d s \leq \int_{0}^{t} 2 e^{-\zeta(t-s)} d s \leq 2 \zeta^{-1}\left(1-e^{-\zeta t}\right), \quad t \geq 0 \\
(-\tilde{w}) r\left(1+\int_{0}^{t}\left|A e^{A(t-s)}\right| d s\right) \leq(-\tilde{w}) r\left(1+2 \zeta^{-1}\left(1-e^{-\zeta t}\right)\right) \leq(-\tilde{w}) r\left(1+2 \zeta^{-1}\right)
\end{gathered}
$$

If $(-\tilde{w}) r / \zeta<1 / 3$, the conditions of Theorem 5.2 are satisfied, that is to say, if $\tilde{w}<0$ very large or $r$ very small, then every solution of (45) and its derivative tends to 0 as $t \rightarrow \infty$.

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[^0]:    Email addresses: guilingmath@gmail.com (Guiling Chen), vangaans@math.leidenuniv.nl (Onno van Gaans), S.M.VerduynLunel@uu.nl (Sjoerd Verduyn Lunel)

