

BIFURCATING POSITIVE STABLE STEADY-STATES FOR A SYSTEM OF DAMPED WAVE EQUATIONS

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Abstract. A system of nonlinear damped wave equations with symmetric linear part is investigated. A positive steady-state bifurcates from the trivial solution as a parameter changes. The spectrum of the linearized operator is studied. Then the stability of the positive steady-state is considered as a solution of the nonlinear hyperbolic system. Asymptotic stability results are found for solutions in R^N , $N \geq 1$. Bifurcation methods are used to find the steady-states, and semigroup methods are used to study stability. Stability results are obtained although the semigroup is not analytic.

1. INTRODUCTION

This article considers the positive steady-state solutions and their stabilities for the following nonlinear system of damped wave equations:

$$\begin{cases} u_{tt} + \beta u_t = \Delta u + \lambda B(x)u + \lambda G(u)u & \text{for } (x,t) \in \Omega \times [0, \infty), \\ u = 0 & \text{for } (x,t) \in \partial\Omega \times [0, \infty). \end{cases} \quad (1.1)$$

Here, $u = \text{col.}(u_1, \dots, u_n)$, $\beta > 0$, λ is a real parameter, $G = [g_{ij}]$ is an $n \times n$ matrix function with each entry in $C^2(R^n)$ and $g_{ij}(0, \dots, 0) = 0$ for $i, j = 1, \dots, n$. Ω is a bounded domain in R^N , with boundary $\partial\Omega$ of class $C^{2+\mu}$, $0 < \mu < 1$. The $n \times n$ matrix $B(x) = [b_{ij}(x)]$ is assumed to satisfy

[H1] $B(x)$ is a real symmetric matrix, with each entry $b_{ij}(x)$ in $C^\mu(\bar{\Omega})$, $0 < \mu < 1$, and nonnegative in $\bar{\Omega}$.

[H2] There is a permutation $\{r_1, r_2, \dots, r_n\}$ of $\{1, 2, \dots, n\}$ such that $b_{r_1 r_2} \not\equiv 0$, $b_{r_2 r_3} \not\equiv 0$, \dots , $b_{r_{n-1} r_n} \not\equiv 0$, and $b_{r_n r_1} \not\equiv 0$ in $\bar{\Omega}$.

(The hypothesis [H1] will be weakened to [H1*] below such that the diagonal entries of $B(x)$ may change sign). Nonlinear hyperbolic systems of similar nature arise in the study of physics and mechanics as the coupled

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sine-Gordon equations (see e.g. [11]). In [10], a recent study is made by energy method concerning the decay of the solutions for similar damped systems. This article can be considered as an extension of the scalar results in [12], where both $n = N = 1$, to systems with symmetric linear part.

In Section 2, we show that a positive steady-state bifurcates from the trivial solution as the parameter λ changes. Moreover, the corresponding elliptic system is linearized at the positive steady-state and the spectrum of the linear operator is investigated under appropriate conditions on the nonlinear terms. In Section 3, the stability of the positive steady-state is considered as a solution of hyperbolic problem (1.1). Note that the corresponding semi-group here is not analytic, and the stability theorem in [4] does not apply. The results in Section 3 concerning stability are shown to be applicable to solutions in R^1 by means of Morrey's inequality, and to solutions in R^3 or R^4 by means of Gagliardo-Nirenberg-Sobolev inequality. In Section 4, positive steady-state is found for $\lambda = 1$ under further assumptions.

For convenience, we will adopt the following conventions. Let $E := \{w = \text{col.}(w_1, \dots, w_n) : w_i \in C^1(\bar{\Omega}), w_i = 0 \text{ on } \partial\Omega, i = 1, \dots, n\}$, $P := \{w = \text{col.}(w_1, \dots, w_n) \in E : w_i \geq 0 \text{ in } \bar{\Omega}, i = 1, \dots, n\}$, and $Y := \{w = \text{col.}(w_1, \dots, w_n) : w_i \in C^{2+\alpha}(\bar{\Omega}), w_i = 0 \text{ on } \partial\Omega, i = 1, \dots, n\}$.

2. BIFURCATION AND SPECTRUM OF LINEARIZED EQUATION

We first show that as the parameter λ passes a certain eigenvalue of the linearized system, a positive steady-state bifurcates from the trivial solution. Then, we will analyze the spectrum of the elliptic part of the operator equation linearized at the positive steady-state.

For convenience, we define the operator L_q with n components as: $L_q \equiv (-\Delta + q_1(x), \dots, -\Delta + q_n(x))$, where $q_i(x), i = 1, \dots, n$ are any nonnegative functions in $C^\mu(\bar{\Omega})$.

Theorem 2.1. *Under hypotheses [H1] and [H2], there exists $(\lambda_0, u^0) \in R \times Y$, such that*

$$L_q[u^0] = \lambda_0 B u^0 \text{ in } \Omega, \quad u^0 = 0 \text{ on } \partial\Omega, \quad (2.1)$$

with $\lambda_0 > 0$, and each component $u_i^0 > 0$ in $\Omega, \partial u_i^0 / \partial \nu < 0$ on $\partial\Omega$ for $i = 1, \dots, n$. Furthermore, $1/\lambda_0$ is a simple eigenvalue of the operator $L_q^{-1}B : [C^1(\bar{\Omega})]^n \rightarrow [C^1(\bar{\Omega})]^n$ (i.e., eigenfunctions corresponding to this eigenvalue is unique up to a multiple). Also, the number $\lambda = \lambda_0$ is the unique positive number so that the problem $u = \lambda L_q^{-1} B u$ has a nontrivial nonnegative solution for $u \in P$.

Proof. The operator $L_q^{-1}B$ is completely continuous and positive with respect to the cone P . By means of Theorem 2.5 in [5] we can obtain a nontrivial $u^0 \in P$ such that $L_q^{-1}Bu^0 = \rho_0 u^0$ for some $\rho_0 > 0$ (i.e., (2.1), with $\lambda_0 = 1/\rho_0$). By the nonnegativity of all entries of B and hypothesis [H2], we show with maximum principle that $u_i^0 > 0$ in Ω and $\partial u_i^0/\partial \nu < 0$ on $\partial\Omega$ for each $i = 1, \dots, n$. By using a comparison principle for systems (Lemma 2.1 in [8]), we can show in the same way as Theorem 2.1 in [8] that $1/\lambda_0$ is a simple eigenvalue, and λ_0 is the unique positive number with a nonnegative eigenfunction as described above.

For a more general situation, we will let the diagonal entries of $B(x)$ to change sign. For $i = 1, \dots, n$, let $b_{ii}(x) = b_{ii}^+(x) + b_{ii}^-(x)$, where $b_{ii}^+(x) = \max\{b_{ii}(x), 0\}$ and $b_{ii}^-(x) = \min\{b_{ii}(x), 0\}$. We introduce following hypothesis:

[H1*] $B(x)$ is a real symmetric matrix with each entry $b_{ij}(x)$ in $C^\mu(\bar{\Omega})$. For all $i \neq j$, b_{ij} is nonnegative in $\bar{\Omega}$, and there exists an integer $k, 1 \leq k \leq n$, such that $b_{kk}^+(x) \not\equiv 0$ in $\bar{\Omega}$.

The following is an extension for Theorem 2.1.

Theorem 2.2. *Under hypotheses [H1*] and [H2], there exists $(\hat{\lambda}_0, v^0) \in R \times Y$ such that*

$$-\Delta v^0 = \hat{\lambda}_0 B(x)v^0 \text{ in } \Omega, \quad v^0 = 0 \text{ on } \partial\Omega, \tag{2.2}$$

with $\hat{\lambda}_0 > 0$, each component $v_i^0 > 0$ in Ω and $\partial v_i^0/\partial \nu < 0$ on $\partial\Omega$ for $i = 1, \dots, n$. Furthermore, $1/\hat{\lambda}_0$ is a simple eigenvalue of the operator $(-\Delta)^{-1}B : [C^1(\bar{\Omega})]^n \rightarrow [C^1(\bar{\Omega})]^n$.

Proof. We will use Theorem 2.1 to prove this theorem. For convenience, define $\tilde{B} = \{\tilde{b}_{ij}(x)\}$ to be the $n \times n$ matrix function on $\bar{\Omega}$ as follows:

$$\tilde{b}_{ij}(x) = b_{ij}(x) \text{ if } i \neq j, \quad \tilde{b}_{ii}(x) = b_{ii}^+(x)$$

for $i, j = 1, \dots, n, x \in \bar{\Omega}$. For each $\lambda \geq 0$, define the n component vector operator

$$\tilde{L}_\lambda \equiv (-\Delta - \lambda b_{11}^-(x), -\Delta - \lambda b_{22}^-(x), \dots, -\Delta - \lambda b_{nn}^-(x))$$

and consider the eigenvalue problem

$$\tilde{L}_\lambda u = \rho \tilde{B}u \text{ in } \Omega, \quad u|_{\partial\Omega} = 0, \tag{2.3}$$

with eigenvalue ρ . Since \tilde{B} satisfies the conditions in Theorem 2.1, the problem (2.3) has a unique positive solution $\rho = \hat{\rho}(\lambda)$, with corresponding eigenfunction u_λ whose components are all positive in Ω .

To proceed with the proof, we need the following two lemmas.

Lemma 2.1. *Under the hypotheses of Theorem 2.2, the function $\hat{\rho}(\lambda)$ is bounded for all $\lambda \in [0, \infty)$.*

Proof. By hypothesis $[H1^*]$, there is an open set D in Ω with its closure in Ω , such that $b_{kk} = b_{kk}^+$ in D for some k . Let Φ be a nontrivial C^∞ function with compact support contained in D . We clearly have $\int_D b_{kk} \Phi^2 dx > 0$. Let u_λ be as described above, and set $w_\lambda(x) = \ln(u_\lambda)_k(x)$ for $x \in \Omega$. Thus, we have in D that

$$\begin{aligned} -\Delta w_\lambda - \sum_{i=1}^N (\partial(w_\lambda)/\partial x_i)^2 &= [1/(u_\lambda)_k] [-\Delta(u_\lambda)_k - b_{kk}^-(x)(u_\lambda)_k] \\ &= [\hat{\rho}(\lambda)/(u_\lambda)_k] \sum_{j=1}^n \tilde{b}_{kj}(x)(u_\lambda)_j \geq b_{kk}(x)\hat{\rho}(\lambda). \end{aligned}$$

Multiplying by Φ^2 and integrating over D by parts on the left, we obtain

$$\int_D \langle \Phi \nabla w_\lambda, 2\nabla \Phi - \Phi \nabla w_\lambda \rangle dx \geq \hat{\rho}(\lambda) \int_D b_{kk} \Phi^2 dx.$$

From this we deduce

$$\left[\int_D \langle \nabla \Phi, \nabla \Phi \rangle dx \right] \left[\int_D b_{kk} \Phi^2 dx \right]^{-1} \geq \hat{\rho}(\lambda) > 0$$

for all $\lambda \in [0, \infty)$.

Lemma 2.2. *Under the hypotheses of Theorem 2.2, the function $\hat{\rho}(\lambda)$ is continuous for $\lambda \in [0, \infty)$.*

The proof of the two lemmas above are similar to that for Lemmas 2.2 and 2.3 in [8]. Further details will thus be omitted.

To complete the proof of Theorem 2.2, we solve the equation $\hat{\rho}(\lambda) - \lambda = 0$ for $\lambda = \hat{\lambda}_0$ as described in (2.2). The solution of the equation must exist due to the properties of $\hat{\rho}(\lambda)$ described in Lemmas 2.1 and 2.2. The simplicity of the eigenvalue $1/\hat{\lambda}_0$ follows from (2.3) and Theorem 2.1.

Note that the pair $(\hat{\lambda}_0, v^0)$ in Theorem 2.2 also satisfies

$$-\Delta v^0 = \hat{\lambda}_0 B^T(x)v^0 \tag{2.4}$$

since $B = B^T$. For convenience, we define an operator $F : R^+ \times E \rightarrow E$ by

$$F(\lambda, u) := u - \lambda(-\Delta)^{-1}[B + G(u)]u \quad \text{for } (\lambda, u) \in R^+ \times E. \tag{2.5}$$

The steady-state solution of (1.1) can be written as

$$F(\lambda, u) = 0. \tag{2.6}$$

Defining operators

$$\begin{aligned} L_0 : E \rightarrow E & \quad \text{by} \quad L_0 := I - \hat{\lambda}_0(-\Delta)^{-1}B, \\ L_1 : E \rightarrow E & \quad \text{by} \quad L_1 := \Delta^{-1}B, \quad \text{and} \\ \tilde{G} : R^+ \times E \rightarrow E & \quad \text{by} \quad \tilde{G}(\lambda, u) := \lambda\Delta^{-1}[G(u)u], \end{aligned}$$

equation (2.6) becomes

$$L_0u + (\lambda - \hat{\lambda}_0)L_1u + \tilde{G}(\lambda, u) = 0, \text{ for } (\lambda, u) \in R^+ \times E. \quad (2.7)$$

As in [8], we readily obtain the following.

Lemma 2.3. *Under the hypotheses of Theorem 2.2, the null space and range of L_0 , denoted respectively by $N(L_0)$ and $R(L_0)$, satisfy:*

- (i) $N(L_0)$ is one-dimensional, spanned by v^0 ;
- (ii) $\dim[E/R(L_0)] = 1$;
- (iii) $L_1v^0 \notin R(L_0)$.

Applying the bifurcation theorem in [1], we obtain the following.

Theorem 2.3. *Assume hypotheses $[H1^*], [H2]$ and each entry of G is in $C^2(R^n)$. Then the point $(\hat{\lambda}_0, 0)$ is a bifurcation point for the problem (2.6). Moreover, there exists a $\delta > 0$ and a C^1 -curve $(\hat{\lambda}(s), \hat{\phi}(s)) : (-\delta, \delta) \rightarrow R \times E$ with $\hat{\lambda}(0) = \hat{\lambda}_0, \hat{\phi}(0) = 0$, such that in a neighborhood of $(\hat{\lambda}_0, 0)$, any solution of (2.6) is either of the form $(\lambda, 0)$ or on the curve $(\hat{\lambda}(s), s[v^0 + \hat{\phi}(s)])$ for $|s| < \delta$, where $s[v^0 + \hat{\phi}(s)] > 0$ in Ω .*

Further applying the theory in [1], we can assert that there exists $\delta_1 \in (0, \delta)$ and two functions $(\sigma(\cdot), z(\cdot)) : (\hat{\lambda}_0 - \delta_1, \hat{\lambda}_0 + \delta_1) \rightarrow R \times E$, and $(\eta(\cdot), h(\cdot)) : [0, \delta_1] \rightarrow R \times E$ with $(\sigma(\hat{\lambda}_0), z(\hat{\lambda}_0)) = (\eta(0), h(0)) = (0, v^0)$ such that

$$\begin{aligned} D_2F(\lambda, 0)z(\lambda) &= \sigma(\lambda)\Delta^{-1}(z(\lambda)), \text{ and} \\ D_2F(\hat{\lambda}(s), s[v^0 + \hat{\phi}(s)])h(s) &= \eta(s)\Delta^{-1}(h(s)). \end{aligned}$$

Here, $\sigma(\lambda)$ and $\eta(s)$ are respectively Δ^{-1} -simple eigenvalues with eigenfunctions $z(\lambda)$ and $h(s)$. The theory in [1] leads to the following.

Lemma 2.4. *Assume all the hypotheses in Theorem 2.3. There exists $\rho > 0$ such that for each $s \in [0, \delta_1)$, there is a unique (real) eigenvalue $\eta(s)$ for the linear operator*

$$F_s^* := \Delta D_2F(\hat{\lambda}(s), s(v^0 + \hat{\phi}(s))) : Y \rightarrow [C^\mu(\bar{\Omega})]^n \quad (2.8)$$

satisfying $|\eta(s)| < \rho$ with eigenfunction $h(s) \in Y$. That is,

$$F_s^*h(s) \equiv \Delta[h(s)] + \hat{\lambda}(s)[B + G(u_s) + G_u(u_s)u_s]h(s) = \eta(s)h(s), \quad (2.9)$$

where $u_s := s(v^0 + \hat{\phi}(s))$. Here, $G_u(u_s)u_s$ denotes the $n \times n$ matrix whose i -th column is $\{\partial G/\partial u_i\}(u_s)u_s$.

In order to determine the signs of the real part of the eigenvalues we need additional assumptions.

Lemma 2.5. *Assume all the hypotheses of Theorem 2.3. Suppose further that*

[H3] $\{\partial g_{ij}/\partial u_k\}(0) \leq 0$ for all $i, j, k = 1, \dots, n$, with at least one inequality being strict. Then the function $\hat{\lambda}(s)$ satisfies $\hat{\lambda}'(0) > 0$.

Proof. Theorem 2.3 asserts that $\hat{\lambda}'(0)$ exists; moreover, for $s \in [0, \delta)$, we have

$$\Delta s(v^0 + \hat{\phi}(s)) + \hat{\lambda}(s)B(x)s(v^0 + \hat{\phi}(s)) + \hat{\lambda}(s)G(s(v^0 + \hat{\phi}(s))s(v^0 + \hat{\phi}(s))) = 0.$$

Dividing by s , then differentiating with respect to s and setting $s = 0$, we obtain

$$\Delta \hat{\phi}'(0) + \hat{\lambda}'(0)B(x)v^0 + \hat{\lambda}_0 B(x)\hat{\phi}'(0) + \hat{\lambda}_0 \frac{d}{ds}G(s(v^0 + \hat{\phi}(s)))|_{s=0} v^0 = 0.$$

Multiplying by $(v^0)^T$ and integrating by parts over Ω , we find

$$\begin{aligned} \hat{\lambda}'(0) &= \frac{-\hat{\lambda}_0 \int_{\Omega} (v^0)^T \frac{d}{ds}G[s(v^0 + \hat{\phi}(s))] |_{s=0} v^0 dx}{\int_{\Omega} (v^0)^T B v^0 dx} \\ &= \frac{-\hat{\lambda}_0 \int_{\Omega} (v^0)^T [\sum_{i=1}^n v_i^0 \frac{\partial G}{\partial u_i}(0)] v^0 dx}{(1/\lambda_0) \int_{\Omega} |\nabla v^0|^2 dx} > 0. \end{aligned}$$

The last inequality is due to hypothesis [H3].

Lemma 2.6. *Under all the hypotheses of Theorem 2.3. The function $\sigma(\lambda)$ satisfies $\sigma'(\hat{\lambda}_0) > 0$.*

The proof of the last lemma is similar to that for Lemma 3.3 in [8], while using the additional fact that $B = B^T$ here. The details will be omitted.

Lemma 2.7. *Under all the hypotheses of Theorem 2.3, and [H3], there exists $\delta_2 \in (0, \delta_1)$ such that $\eta(s) < 0$ for all $s \in (0, \delta_2)$.*

Proof. From Theorem 1.16 in [1], we find $-s\hat{\lambda}'(s)\sigma'(\hat{\lambda}_0)$ and $\eta(s)$ have the same sign for $s > 0$ near 0. The conclusion follows from Lemmas 2.5 and 2.6.

The linearized eigenvalue problem for (2.6) at the bifurcating solution $u = s(v^0 + \hat{\phi}(s))$ is precisely (2.9). When $s = 0$, $\lambda = \hat{\lambda}(0) = \hat{\lambda}_0$, the eigenvalue problem corresponding to (2.9) becomes

$$\Delta[h] + \hat{\lambda}_0 B h = \eta h \quad h \in Y \quad (2.10)$$

where η is the eigenvalue. Under hypotheses $[H1^*]$ and $[H2]$, Theorem 2.2 asserts that $\eta = 0$ is an eigenvalue for (2.10) with positive eigenfunction. Using this property and the fact that the off-diagonal terms of B are all nonnegative, we can show the following as in [7].

Lemma 2.8. *Under hypotheses $[H1^*]$ and $[H2]$, all eigenvalues in equation (2.10) except $\eta = 0$ satisfies $Re(\eta) < -r$ for some positive number r .*

Theorem 2.4. *Under the hypotheses of Theorem 2.3 and $[H3]$, there exists a number $\delta^* \in (0, \delta)$ and a positive function $\hat{\eta}(s)$ for $s \in (0, \delta^*)$ such that the real parts of all the numbers in the point spectrum of the linear operator F_s^* are contained in the interval $(-\infty, -\hat{\eta}(s))$, for $s \in (0, \delta^*)$. (Here, δ is described in Theorem 2.3 and F_s^* is described in (2.8) in Lemma 2.4).*

The proof of Lemma 2.8 is the same as the proof of Lemma 2.8 in [7]; and the proof of Theorem 2.4, using the assertions in Lemma 2.4 to 2.8 and perturbation arguments, is the same as Theorem 2.2 in [7]. The details are thus omitted. More details of proof of lemmas and theorems similar to those in this section can also be found in [8].

3. STABILITY OF POSITIVE STEADY-STATES

For each $s \in (0, \delta^*)$, the function $u_s := s[v^0 + \hat{\phi}(s)]$ described in Theorem 2.3 and Lemma 2.4 can be considered as a steady-state solution of (1.1) with $\lambda = \hat{\lambda}(s)$. We now consider the time asymptotic stability of this steady state as a solution of the system of hyperbolic equations (1.1). We convert (1.1) into a first order system by letting

$$\begin{cases} \partial u_i / \partial t = v_i & i = 1, \dots, n, \\ \partial v_i / \partial t = \Delta u_i - \beta v_i + \lambda \sum_{j=1}^n [b_{ij}(x) + g_{ij}(u_1, \dots, u_n)] u_j. \end{cases} \tag{3.1}$$

Let $J(u_s(x))$ be the linearization of $G(u)u$ (i.e., $\partial[G(u)u]/\partial u$) evaluated at $u = u_s$. For convenience, define $\bar{B}(u_s(x)) := B(x) + J(u_s(x))$ and let \bar{b}_{ij} denotes the ij -th entry of $\bar{B}(u_s(x))$. The system (3.1) linearized at u_s and $\lambda = \hat{\lambda}(s)$ can be written in the form

$$\partial \xi / \partial t = A_s \xi, \tag{3.2}$$

where $\xi = \text{col.}(\hat{u}_1, \hat{v}_1, \dots, \hat{u}_n, \hat{v}_n)$, and A_s is the differential linear operator on the $2n$ components of ξ as follows:

$$A_s = \begin{bmatrix} A_{11}^s & \dots & A_{1n}^s \\ \dots & \dots & \dots \\ A_{n1}^s & \dots & A_{nn}^s \end{bmatrix}, \text{ where } A_{ij}^s \text{ are } 2 \times 2 \text{ blocks}$$

$$A_{ii}^s = \begin{bmatrix} 0 & 1 \\ \Delta + \hat{\lambda}(s)\bar{b}_{ii} & -\beta \end{bmatrix}, \quad A_{ij}^s = \begin{bmatrix} 0 & 0 \\ \hat{\lambda}(s)\bar{b}_{ij} & 0 \end{bmatrix} \text{ for } i \neq j, \quad i, j = 1, \dots, n.$$

Here, each $A_{i_j}^s$ can be considered as an operator from $[H^2(\Omega) \cap H_0^1(\Omega)] \times H_0^1(\Omega)$ into $H_0^1(\Omega) \times L^2(\Omega)$. Direct calculation shows that if θ is an eigenvalue of A_s , then $\theta^2 + \beta\theta$ is an eigenvalue of F_s^* . That is, if $A_s \bar{\xi} = \theta \bar{\xi}$ for some $\bar{\xi} \in ([H^2(\Omega) \cap H_0^1(\Omega)] \times H_0^1(\Omega))^n$, then

$$\Delta \bar{u} + \hat{\lambda}(s) \bar{B}(u_s) \bar{u} = (\theta^2 + \beta\theta) \bar{u},$$

where $\bar{\xi} = \text{col}(\bar{u}_1, \bar{v}_1, \dots, \bar{u}_n, \bar{v}_n)$ and $\bar{u} = \text{col}(\bar{u}_1, \dots, \bar{u}_n)$. In order to obtain the stability of the steady-state u_s , we will impose an additional assumption:

[H4] $\bar{B}(u_s(x))$ is symmetric in $\bar{\Omega}$ for each $s \in (0, \delta^*)$.

By Theorem 2.4 and hypotheses [H4], as an operator on Y , the eigenvalues of $F_s^* = \Delta + \hat{\lambda}(s) \bar{B}(u_s)$ are all strictly negative for $s \in (0, \delta^*)$. However, the eigenvalues of $\Delta + \hat{\lambda}(s) \bar{B}(u_s)$ are the same, as an operator on Y or $[H^2(\Omega) \cap H_0^1(\Omega)]^n$. Consequently, the eigenvalues θ of A_s , $s \in (0, \delta^*)$, satisfies

$$\theta^2 + \beta\theta - \eta = 0,$$

where η are negative real numbers. Thus, the real parts of θ are all strictly negative. The operator A_s generates a strongly continuous semigroup $T_s(t)$ on $X := [H_0^1(\Omega) \times L^2(\Omega)]^n$ with domain $D(A_s) := ([H^2(\Omega) \cap H_0^1(\Omega)] \times H_0^1(\Omega))^n$. Using the fact that the real parts of θ are negative, we now show the following stability result for the linearized system (3.2).

Theorem 3.1. *Assume hypotheses [H1*], [H2], [H3] and each entry of G is in $C^2(\mathbb{R}^n)$. Suppose further that $\bar{B}(u_s(x))$ is symmetric as described in [H4], then the semigroup of bounded linear operator $T_s(t)$, $0 \leq t < \infty$, on X satisfies*

$$\|T_s(t)\| \leq M e^{-kt}$$

for some positive constants M and k , which may depend on s .

Proof. Let p be a large enough positive constant such that

$$y^T [\hat{\lambda}(s) \bar{B}(u_s(x)) - pI] y \leq -\omega (y^T \cdot y)$$

for all $s \in (0, \delta^*)$, $x \in \bar{\Omega}$, where ω is a positive constant. Define the operator $S := -(\Delta + \hat{\lambda} \bar{B} - pI)^{-1}$ from $[L^2(\Omega)]^n$ into $[L^2(\Omega)]^n$ as follows. For each $f \in [L^2(\Omega)]^n$, let $g := S[f]$, where $g \in [H_0^1(\Omega)]^n \cap [H^2(\Omega)]^n$, and

$$\langle\langle \nabla g, \nabla h \rangle\rangle + \int_{\Omega} h^T (-\hat{\lambda} \bar{B} + pI) g dx = \int_{\Omega} h^T f dx$$

for all $h \in [H_0^1(\Omega)]^n$. Here $\langle\langle \cdot \cdot \rangle\rangle$ is the real inner product in $[L^2(\Omega)]^{nN}$. The operator $S : [L^2(\Omega)]^n \rightarrow [L^2(\Omega)]^n$ is well-defined and compact by means of the Lax-Milgram Theorem and Sobolev imbedding (cf. [3]). From the

symmetric property of \bar{B} , we can readily verify that the operator S is symmetric as follows. Let $f, h \in [L^2(\Omega)]^n$, $g = S[f], q = S[h]$, and $\langle \cdot \rangle$ denotes the real inner product in $[L^2(\Omega)]^n$. Then we have $\langle Sf, h \rangle = \langle g, -(\Delta + \hat{\lambda}\bar{B} - pI)q \rangle = \langle -\Delta g, q \rangle + \langle g, -(\hat{\lambda}\bar{B} - pI)q \rangle = \langle -\Delta g, q \rangle + \langle -(\hat{\lambda}\bar{B} - pI)g, q \rangle = \langle f, q \rangle = \langle f, Sh \rangle$. From the theory of compact symmetric operators, we assert that there exists a countable orthonormal basis of $[L^2(\Omega)]^n$ consisting of eigenfunctions of S . Direct computation shows that if η is an eigenvalue of $(\Delta + \hat{\lambda}\bar{B})$, then $1/(p - \eta)$ is an eigenvalue of S . Theorem 2.4 and the additional assumption that \bar{B} is symmetric thus imply that all the eigenvalues of S are of the form $1/(p + \alpha_i)$, with $0 < \alpha_1 < \alpha_2 < \dots$, (Note that α_i depends on s). There exist corresponding eigenfunctions $\{\phi_m\}_{m=1}^\infty$, which form an orthonormal basis in $[L^2(\Omega)]^n$.

Let $[\tilde{H}_0^1(\Omega)]^n$ denotes the real Hilbert space of functions in $[H_0^1(\Omega)]^n$ with inner product

$$\tilde{J}[g, q] := \langle \nabla g, \nabla q \rangle + \int_{\Omega} q^T [-\hat{\lambda}\bar{B} + pI]g \, dx.$$

We verify that

$$\tilde{J}[\phi_j, \phi_k] = \langle -\Delta\phi_j, \phi_k \rangle + \int_{\Omega} \phi_k^T [-\hat{\lambda}\bar{B} + pI]\phi_j \, dx = \langle (\alpha_j + p)\phi_j, \phi_k \rangle.$$

Thus, $\{\phi_m/(\alpha_m + p)^{1/2}\}_{m=1}^\infty$ form an orthonormal set in $[\tilde{H}_0^1(\Omega)]^n$. Further, the identity

$$\tilde{J}[\phi_j, g] = \int_{\Omega} g^T [-\Delta\phi_j - \hat{\lambda}\bar{B} + pI]\phi_j \, dx = \langle (\alpha_j + p)\phi_j, g \rangle$$

implies that if $\tilde{J}[\phi_j, g] = 0$ for each $j = 1, 2, \dots$, then we also have $\langle \phi_j, g \rangle = 0$. This implies that $\{\phi_m/(\alpha_m + p)^{1/2}\}_{m=1}^\infty$ form an orthonormal basis in $[\tilde{H}_0^1(\Omega)]^n$. If $\hat{u} \in [\tilde{H}_0^1(\Omega)]^n$, we can assert that the series

$$\hat{u} = \sum_{m=1}^{\infty} \langle \hat{u}, \phi_m \rangle \phi_m \quad \text{converges in } [L^2(\Omega)]^n,$$

and further

$$\begin{aligned} \hat{u} &= \sum_{m=1}^{\infty} \tilde{J}[\hat{u}, \phi_m/(\alpha_m + p)^{1/2}]\{\phi_m/(\alpha_m + p)^{1/2}\} \\ &= \sum_{m=1}^{\infty} (\alpha_m + p)^{1/2} \langle \hat{u}, \phi_m \rangle \{\phi_m/(\alpha_m + p)^{1/2}\} = \sum_{m=1}^{\infty} \langle \hat{u}, \phi_m \rangle \phi_m \end{aligned}$$

actually converges in $[\tilde{H}_0^1(\Omega)]^n$.

For $\phi \in D(A_s)$, let $\hat{w} = \hat{w}(t; \phi) := T_s(t)\phi \in D(A_s)$, we know from semigroup theory that $d\hat{w}/dt \in X$ for $t > 0$. That is, if we let $\hat{w} = (\hat{u}_1, \hat{v}_1, \dots, \hat{u}_n, \hat{v}_n)$, and consider \hat{u}_i as functions from $[0, \infty)$ into X , then $d\hat{u}_i/dt \in H_1^0(\Omega)$, and $d^2\hat{u}_i/dt^2 \in L^2$. This implies that if we let $\hat{u} = \text{col.}(\hat{u}_1, \dots, \hat{u}_n)$, $\hat{v} = \text{col.}(\hat{v}_1, \dots, \hat{v}_n)$ and $r_m(t) := \langle \hat{u}(t), \phi_m \rangle$, then the series

$$d\hat{u}/dt = \sum_{m=1}^{\infty} r'_m(t)\phi_m \quad \text{converges in } [H_1^0(\Omega)]^n, \text{ and}$$

$$d^2\hat{u}/dt^2 = \sum_{m=1}^{\infty} r''_m(t)\phi_m \quad \text{converges in } [L^2(\Omega)]^n,$$

for $t > 0$, where $r'_m(t)$ and $r''_m(t)$ are the first and second derivatives of $r_m(t)$. From semigroup theory and the structure of the operator A_s , we find that $\hat{u}(t)$ satisfies

$$d^2\hat{u}/dt^2 + \beta d\hat{u}/dt = [\Delta + \hat{\lambda}(s)\bar{B}(u_s)]\hat{u}, \quad \text{for } t > 0.$$

Taking inner product with ϕ_m , we obtain

$$r''_m(t) + \beta r'_m(t) = -\alpha_m r_m(t) \quad \text{for } t > 0, m = 1, 2, \dots \quad (3.3)$$

Denote $\phi := \text{col.}(g_1, h_1, \dots, g_n, h_n) \in D(A_s) = [(H^2 \cap H_0^1) \times H_0^1]^n$, with

$$\begin{aligned} g &= \sum_{m=1}^{\infty} \langle g, \phi_m \rangle \phi_m \\ &= \sum_{m=1}^{\infty} (\alpha_m + p)^{1/2} \langle g, \phi_m \rangle [\phi_m / (\alpha_m + p)^{1/2}] \quad \text{in } [H_0^1(\Omega)]^n, \\ h &= \sum_{m=1}^{\infty} \langle h, \phi_m \rangle \phi_m \quad \text{in } [L^2(\Omega)]^n. \end{aligned}$$

Let $T_s(t)\phi := \text{col.}(\hat{u}_1, \hat{v}_1, \dots, \hat{u}_n, \hat{v}_n)$. From (3.3), we find

$$\begin{aligned} r_m(t) &= \langle g, \phi_m \rangle y_m(t) + \langle h, \phi_m \rangle z_m(t) \\ r'_m(t) &= \langle g, \phi_m \rangle y'_m(t) + \langle h, \phi_m \rangle z'_m(t) \end{aligned}$$

where $y_m(t)$ and $z_m(t)$ satisfy the same equations for $r_m(t)$ in (3.3), with initial conditions:

$$y_m(0) = 1, \quad y'_m(0) = 0 \quad \text{and} \quad z_m(0) = 0, \quad z'_m(0) = 1.$$

For those m with $4\alpha_m > \beta^2$, we have

$$\begin{aligned} y_m(t) &= e^{-\beta t/2} \cos(4\alpha_m - \beta^2)^{1/2} t/2, \\ z_m(t) &= 2(4\alpha_m - \beta^2)^{-1/2} e^{-\beta t/2} \sin(4\alpha_m - \beta^2)^{1/2} t/2. \end{aligned} \tag{3.4}$$

Expanding \hat{u} in $[\tilde{H}_0^1(\Omega)]^n$, and using (3.4) for all large m , we find

$$\begin{aligned} \|\hat{u}\|_{\tilde{H}^1} &= \sum_{m=1}^{\infty} [\langle g, \phi_m \rangle y_m(t) + \langle h, \phi_m \rangle z_m(t)]^2 (\alpha_m + p) \\ &\leq \hat{K} e^{-\epsilon t} \left\{ \sum_{m=1}^{\infty} \langle g, \phi_m \rangle^2 (\alpha_m + p) + \sum_{m=1}^{\infty} \langle h, \phi_m \rangle^2 \right\}, \\ \|\hat{v}\|_{L^2} &= \sum_{m=1}^{\infty} [\langle g, \phi_m \rangle y'_m(t) + \langle h, \phi_m \rangle z'_m(t)]^2 \\ &\leq \hat{K} e^{-\epsilon t} \left\{ \sum_{m=1}^{\infty} \langle g, \phi_m \rangle^2 (\alpha_m + p) + \sum_{m=1}^{\infty} \langle h, \phi_m \rangle^2 \right\}. \end{aligned}$$

Since $[\tilde{H}_0^1(\Omega)]^n$ and $[H_0^1(\Omega)]^n$ are equivalent, we obtain

$$\|T_s(t)\phi\|_X \leq M e^{-\epsilon t} \|\phi\|_X$$

for all ϕ in $D(A_s)$. Since $D(A_s)$ is dense in X , the proof of Theorem 3.1 is complete.

If $(u_1(t), v_1(t), \dots, u_n(t), v_n(t))$ is a solution of (3.1), then its difference with the steady-state, i.e., $\text{col}(w_1(t), w_2(t), \dots, w_{2n}(t)) := \text{col}(u_1(t), v_1(t), \dots, u_n(t), v_n(t)) - \text{col}((u_s)_1, 0, \dots, (u_s)_n, 0)$, can be interpreted as a solution of:

$$\frac{dw}{dt} = A_s w + \hat{\lambda}(s)\gamma(w), \quad w = \text{col}(w_1, w_2, \dots, w_{2n}), \tag{3.5}$$

where

$$\begin{aligned} \gamma(w) &:= \text{col}(0, z_1(\tilde{w}), 0, z_2(\tilde{w}), \dots, 0, z_n(\tilde{w})), \\ \tilde{w} &= \text{col}(\tilde{w}_1, \tilde{w}_2, \dots, \tilde{w}_n) := \text{col}(w_1, w_3, \dots, w_{2n-1}), \\ z &= \text{col}(z_1(\tilde{w}), z_2(\tilde{w}), \dots, z_n(\tilde{w})) \\ &:= [G(u_s + \tilde{w})(u_s + \tilde{w}) - G(u_s)u_s - J(u_s)\tilde{w}]. \end{aligned}$$

We now clarify some terminologies and hypotheses which we will be using. A function $h : X \rightarrow X$ is said to satisfy a local Lipschitz condition if for every constant $c \geq 0$, there is a constant $L(c)$ such that

$$\|h(w^*) - h(w^{**})\|_X \leq L(c) \|w^* - w^{**}\|_X$$

holds for all $w^*, w^{**} \in X$ with $\|w^*\|_X \leq c$ and $\|w^{**}\|_X \leq c$. The function $\gamma(w)$ described above is said to satisfy property $[P]$ if

$$[P] \quad \|z(\tilde{w})\|_{[L^2(\Omega)]^n} = o(\|\tilde{w}\|_{[H^1(\Omega)]^n}), \text{ as } \|\tilde{w}\|_{[H^1(\Omega)]^n} \text{ tends to } 0.$$

Remark 3.1. Since the norm in X uses the L^2 -norm of the even components and the entries of $G(u)$ is in $C^2(R^n)$, it can be readily verified by using Sobolev imbedding or Morrey’s inequality that the function $\gamma : X \rightarrow X$ satisfies a local Lipschitz condition, in the case $N = 1$. By Theorem 1.4 of Chapter 6 in [9], equation (3.5) has a unique mild solution w on $[0, t_{\max})$, for every given initial condition: $w(0) = w_0 \in X$. Moreover, if $t_{\max} < \infty$, then $\|w(t)\|_X \rightarrow \infty$, as $t \rightarrow \infty$.

Remark 3.2. Using the assumption that the entries of G are in C^2 and Sobolev imbedding again, we can readily verify that in the case $N = 1$, the function $\gamma(w)$ described above satisfies property $[P]$.

Remark 3.3. The even components of $\gamma(w)$ are expressed by $z(\tilde{w})$, where \tilde{w} is an n -vector function consisting of the odd components of w . Careful calculations shows that the j -th component, $1 \leq j \leq n$, of $z(\tilde{w})$ can be written as:

$$\begin{aligned} & \sum_{k=1}^n (u_s)_k \langle \nabla g_{jk}(u_s + \tau_{jk}^* \tilde{w}) - \nabla g_{jk}(u_s), \tilde{w} \rangle \\ & + \sum_{k=1}^n \tilde{w}_k (g_{jk}(u_s + \tilde{w}) - g_{jk}(u_s)) \\ & = \sum_{k=1}^n (u_s)_k \sum_{i=1}^n \langle (\nabla(\partial g_{jk}/\partial u_i))(u_s + \tau_{jk}^{**} \tilde{w}), \tau_{jk}^* \tilde{w} \rangle \tilde{w}_i \\ & + \sum_{k=1}^n \tilde{w}_k \langle (\nabla g_{jk})(u_s + \hat{\tau}_{jk} \tilde{w}), \tilde{w} \rangle, \end{aligned}$$

where $0 < \tau_{jk}^{**} < \tau_{jk}^* < 1, 0 < \hat{\tau}_{jk} < 1$. When the space dimension N is 3 or 4, Gagliardo-Nirenberg-Sobolev inequality asserts that the L^4 norm of a function in $\Omega \subset R^N$ is bounded by a constant multiple of its $H^1(\Omega)$ norm. If we assume all the first and second partial derivatives of $g_{jk}, 1 \leq j, k \leq n$, are bounded in R^n , then from the above formula, we can readily obtain

$$\|\gamma(w)\|_X = \|z(\tilde{w})\|_{[L^2(\Omega)]^n} \leq K(\|\tilde{w}\|_{[L^4(\Omega)]^n})^2$$

for all $w \in X = [H_0^1(\Omega) \times L^2(\Omega)]^n$. Thus, property $[P]$ is satisfied. The local Lipschitz condition of $\gamma : X \rightarrow X$ can be verified similarly. In short, if all the entries of $G(u)$ are in $C^2(R^n)$ and has bounded first and second partial

derivatives in R^n , then the function γ satisfies the local Lipschitz condition and property $[P]$ described above, provided the space dimension N of Ω is 3 or 4.

Using Gronwall inequality type argument, we then use semigroup theory as in [4] or [9] to prove the following asymptotic stability theorem by means of Theorem 3.1.

Theorem 3.2. *Assume hypotheses $[H1^*], [H2], [H3], G$ is in $C^2(R^n)$, and $\bar{B}(u_s)$ is symmetric for each $s \in (0, \delta^*)$ as described in $[H4]$. Suppose that in equation (3.5), the function γ satisfies a local Lipschitz condition and property $[P]$, then the trivial steady-state solution $w = 0$ of equation (3.5) is locally asymptotically stable. That is, there exist positive constants $\epsilon, \hat{M}, \hat{k}$ such that if the initial condition $w(0) := (w_1(0), w_2(0), \dots, w_{2n}(0))$ satisfies $\|w(0)\|_X \leq \epsilon$, then the unique mild solution of (3.5) satisfies*

$$\|w(t)\|_X \leq \hat{M}e^{-\hat{k}t}\|w(0)\|_X. \tag{3.6}$$

Proof. The local Lipschitz condition of γ insures the existence of a unique mild solution $w(t) \in X$ of (3.5) as described in Theorem 1.4 in Chapter 6 in [9]. The mild solution $w(t)$ can be expressed by:

$$w(t) = T_s(t)w(0) + \int_0^t T_s(t-\tau)\hat{\lambda}(s)\gamma(w(\tau))d\tau, \quad w(0) \in X, \tag{3.7}$$

where T_s is the strongly continuous semigroup generated by A_s . Let M and k be the positive constants given by Theorem 3.1 for T_s , we obtain for $t > 0$ inside the interval of existence of $w(t)$:

$$\|w(t)\|_X \leq Me^{-kt}\|w(0)\|_X + \int_0^t Me^{-k(t-\tau)}\|\hat{\lambda}(s)\gamma(w(\tau))\|_Xd\tau. \tag{3.8}$$

Without loss of generality, we may assume $M > 1$. By property $[P]$ for γ , there exist $\rho > 0$ so small such that

$$\|\hat{\lambda}(s)\gamma(w)\|_X \leq [k/(2M)]\|\tilde{w}\|_{[H^1]^n} \leq [k/(2M)]\|w\|_X \tag{3.9}$$

as long as $\|w\|_X \leq \rho$. Let $\|w(0)\|_X \leq \rho/M$, the solution will satisfy $\|w(t)\|_X \leq \rho$ for sufficiently small $t > 0$. For such small $t > 0$, (3.8) and (3.9) lead to

$$e^{kt}\|w(t)\|_X \leq M\|w(0)\|_X + (k/2) \int_0^t e^{k\tau}\|w(\tau)\|_Xd\tau. \tag{3.10}$$

Thus, Gronwall's inequality gives

$$\|w(t)\|_X \leq M\|w(0)\|_Xe^{-(k/2)t} \leq \rho e^{-(k/2)t} \tag{3.11}$$

as long as $\|w(t)\|_X \leq \rho$. Thus, $w(t)$ exists for all $t > 0$, if we choose $w(0)$ to satisfy $\|w(0)\|_X \leq \rho/M$; and (3.11) is valid for all $t > 0$ in such cases. Consequently, Theorem 3.2 is valid by choosing $\epsilon = \rho/M$, $\hat{M} = M$, and $\hat{k} = k/2$.

We now return to equation (3.1), which is the first order system form of (1.1). Equation (3.1) can be written as:

$$\frac{\partial y}{\partial t} = \tilde{A}_s y + \hat{\lambda}(s)q(y), \quad (3.12)$$

where $y = \text{col}(y_1, y_2, \dots, y_{2n}) = \text{col}(u_1, v_1, \dots, u_n, v_n)$, \tilde{A}_s is the same as A_s with all \bar{b}_{ij} replaced by the corresponding b_{ij} , and $q(y) := \text{col}(0, \hat{q}_1(\hat{y}), \dots, 0, \hat{q}_n(\hat{y}))$, $\hat{y} = \text{col}(y_1, y_3, \dots, y_{2n-1})$, $\text{col}(\hat{q}_1(\hat{y}), \dots, \hat{q}_n(\hat{y})) := G(\hat{y})\hat{y}$. We will express the results in Theorem 3.2 in terms of the mild solutions in X of the semigroup form of (3.12), that is,

$$\frac{dy}{dt} = \tilde{A}_s y + \hat{\lambda}(s)q(y), \quad y(t) \in X, t \geq 0. \quad (3.13)$$

The constant function $\eta(t) := \text{col}((u_s)_1, 0, \dots, (u_s)_n, 0)$ is a strong solution of

$$\frac{d\eta}{dt} = A_s \eta + \hat{\lambda}(s)h(u_s), \quad (3.14)$$

where $h(u_s) = \text{col}(0, \hat{h}_1(u_s), \dots, 0, \hat{h}_n(u_s))$, and $\text{col}(\hat{h}_1(u_s), \dots, \hat{h}_n(u_s)) = G(u_s)u_s - J(u_s)u_s$. Thus, $\eta(t)$ satisfies

$$\eta(t) = T_s(t)\eta(0) + \int_0^t T_s(t-\tau)\hat{\lambda}(s)h(u_s)d\tau.$$

If $w(t)$ is a mild solution of (3.5), then $w(t) + \eta$ satisfies

$$w(t) + \eta = T_s(t)[w(0) + \eta] + \int_0^t T_s(t-\tau)\hat{\lambda}(s)\{\gamma(w(\tau)) + h(u_s)\}d\tau. \quad (3.15)$$

If $w(0) \in D(A_s)$, then $y = w(t) + \eta$ is a strong solution of

$$\frac{dy}{dt} = A_s y + \hat{\lambda}(s)[\gamma(w) + h(u_s)] = \tilde{A}_s y + \hat{\lambda}(s)[\gamma(w) + h(u_s) + r(y)],$$

where $r(y) = \text{col}(0, \hat{r}_1(\hat{y}), \dots, 0, \hat{r}_n(\hat{y}))$, $\text{col}(\hat{r}_1(\hat{y}), \dots, \hat{r}_n(\hat{y})) = J(u_s)\hat{y}$. Thus, if $w(0) \in D(A_s)$, we also have

$$\begin{aligned} w(t) + \eta &= \tilde{T}_s(t)[w(0) + \eta] \\ &+ \int_0^t \tilde{T}_s(t-\tau)\hat{\lambda}(s)\{\gamma(w(\tau)) + h(u_s) + r(w(\tau) + \eta)\}d\tau, \end{aligned} \quad (3.16)$$

where \tilde{T}_s is continuous semigroup generated by \tilde{A}_s . By the density of $D(A_s)$ in X and the strong continuity of T_s and \tilde{T}_s , the right-hand sides of both (3.15) and (3.16) are still equal for all cases with $w(0) \in X$. Since $\gamma(w) + h(u_s) + r(w + \eta) = q(w + \eta)$, we obtain from (3.16) that $w(t) + \eta$ is a mild solution of equation (3.13).

Theorem 3.2 and the above arguments lead to the following conclusion concerning solutions of (1.1) or (3.1).

Theorem 3.3. *Assume all the hypotheses of Theorem 3.2 (including those concerning γ). Then if the initial condition*

$$(y_1(0), y_2(0), \dots, y_{2n}(0)) := (u_1(0), v_1(0), \dots, u_n(0), v_n(0))$$

is given sufficiently close to $((u_s)_1, 0, \dots, (u_s)_n, 0)$ in X , the unique mild solution $y(t) = (u_1(t), v_1(t), \dots, u_n(t), v_n(t))$ of the initial value problem corresponding to (3.13) exists for all $t > 0$. Moreover, there exist positive constants $\epsilon, \hat{M}, \hat{k}$ such that if

$$\|(u_1(0), v_1(0), \dots, u_n(0), v_n(0)) - ((u_s)_1, 0, \dots, (u_s)_n, 0)\|_X \leq \epsilon,$$

then the mild solution of (3.13) satisfies

$$\|(u_1(t), v_1(t), \dots, u_n(t), v_n(t)) - ((u_s)_1, 0, \dots, (u_s)_n, 0)\|_X \leq$$

$$\hat{M}e^{-\hat{k}t} \|((u_s)_1, 0, \dots, (u_s)_n, 0)\|_X$$

for all $t > 0$. Recall that (3.13) is the semigroup form of (1.1) or (3.1).

For applications of the theorems above, consider the following examples.

Example 3.1. Let $\Omega = \{x = (x_1, x_2, x_3) : x_1^2 + x_2^2 + x_3^2 < 1\}$. Consider (1.1), with $\beta = 1, u = \text{col}(u_1, u_2, u_3)$,

$$B(x) = \begin{bmatrix} \sin 2\pi(x_1^2 + x_2^2 + x_3^2) & 1 & 2 \\ 1 & \cos 2\pi(x_1^2 + x_2^2 + x_3^2) & 3 \\ 2 & 3 & -(x_1^2 + x_2^2 + x_3^2) \end{bmatrix},$$

$$G(u)u = \begin{bmatrix} -u_1 & -u_1 & 0 \\ -u_1/2 & -\sin u_2 & 0 \\ 0 & 0 & 0 \end{bmatrix} \begin{bmatrix} u_1 \\ u_2 \\ u_3 \end{bmatrix} = \begin{bmatrix} -u_1^2 - u_1u_2 \\ -(1/2)u_1^2 - u_2 \sin u_2 \\ 0 \end{bmatrix}.$$

It is clear that $B(x)$ satisfies $[H1^*]$ and $[H2]$, $G(u)$ satisfies $[H3]$ with every entry in C^2 , and

$$\bar{B}(u_s(x)) = B(x) + \begin{bmatrix} -2(u_s)_1 - (u_s)_2 & -(u_s)_1 & 0 \\ -(u_s)_1 & -(u_s)_2 \cos(u_s)_2 - \sin(u_s)_2 & 0 \\ 0 & 0 & 0 \end{bmatrix}$$

is symmetric as specified in [H4]. Since the first and second partial derivatives of $G(u)$ are bounded, by Remark 3.3 the function γ satisfies the local Lipschitz condition and property [P]. Thus, all assumptions of Theorem 3.3 are satisfied, and the steady-state $((u_s)_1, 0, (u_s)_2, 0, (u_s)_3, 0)$ is locally asymptotically stable in the sense described by Theorem 3.3 for the equation (3.13), which is the semigroup form for (1.1). Here, we are considering those $\lambda = \hat{\lambda}(s)$ close to the right of $\hat{\lambda}_0$, $s \in (0, \delta^*)$, as described in Theorem 2.4.

Example 3.2. Let $\Omega = [0, \pi]$, consider (1.1) with $\beta = 1, u = \text{col}(u_1, u_2, u_3)$,

$$B(x) = \begin{bmatrix} 5 & 6 & 7 \\ 6 & 10 & 4 \\ 7 & 4 & -1 \end{bmatrix}$$

$$G(u)u = \begin{bmatrix} -u_1 & -\sin(u_1u_2) & 0 \\ -\sin(u_1u_2) & -u_2^2 & 0 \\ 0 & 0 & 0 \end{bmatrix} \begin{bmatrix} u_1 \\ u_2 \\ u_3 \end{bmatrix} = \begin{bmatrix} -u_1^2 - u_2 \sin(u_1u_2) \\ -u_1 \sin(u_1u_2) - u_2^3 \\ 0 \end{bmatrix}.$$

It can be verified readily that $\bar{B}(u_s(x))$ is symmetric and all the hypotheses of Theorem 3.3 are satisfied. Here, $N = 1$, by Remarks 3.1 and 3.2 we find γ satisfies the local Lipschitz condition and [P], without requiring the boundedness of the first and second partial derivatives of G .

4. GLOBAL BIFURCATION AND EXAMPLE

In this section, we consider the global behavior of the bifurcation curve of positive solution of equation (2.6). We find conditions so that equation (1.1) has a positive steady-state solution with $\lambda = 1$. We now introduce a few notation and hypotheses. Let $\mathcal{G} := \{(\lambda, w) \in R^+ \times P : F(\lambda, w) = 0, \lambda > 0 \text{ and } w \in P \setminus \{0\}\}$, $\bar{\mathcal{G}}$ denotes the closure of \mathcal{G} ; and define the following conditions:

[A1] There exists an integer $m, 1 \leq m \leq n$ such that $b_{mj}(x) > 0, b_{jm}(x) > 0$ in $\bar{\Omega}$ for all $j \neq m$, and the m -th component of $G(u)u$ is $\equiv 0$ in $\bar{\Omega}$.

[A2] For each $j \neq m$, the j -th component of $G(u)u$ is expressible as $k_j(u_1, \dots, u_n)u_j$ with $k_j(u_1, \dots, u_n) \in C^0$ for $u_i \geq 0, i = 1, \dots, n$.

Theorem 4.1. *Under hypotheses [H1], [H1*], [H2], [A1], and [A2], the component of $\bar{\mathcal{G}}$ containing the point $(\lambda_0, 0)$ is unbounded, and $\bar{\mathcal{G}} \cap (R \times \partial P) = (\lambda_0, 0)$.*

Proof. Theorem 2.1 and (2.5) imply by means of Theorem 29.2 in [2] that the component of $\bar{\mathcal{G}}$ containing the point $(\lambda_0, 0)$ is unbounded. Let $(\lambda_i, w_i) \in \mathcal{G}, i = 1, 2, \dots$ be a sequence tending to a limit point $(\bar{\lambda}, \bar{w})$ in $R \times \partial P$, and $(\bar{\lambda}, \bar{w}) \neq (\lambda_0, 0)$. We now show $\bar{w} = \text{col}(\bar{w}_1, \dots, \bar{w}_n)$ must

satisfy $\bar{w}_i \equiv 0$, for each i , if $\bar{\lambda} > 0$. Consider the first case when there exists some $x_0 \in \Omega$, where $\bar{w}_m(x_0) = 0$. The equation

$$-\Delta \bar{w}_m = \bar{\lambda} \sum_{j=1}^n b_{mj} \bar{w}_j \geq 0$$

implies that $\bar{w}_m \equiv 0$ in $\bar{\Omega}$; and subsequently, the right hand side of this equation and [A1] imply that $\bar{w}_j \equiv 0$ for each $j \neq m$ too. Hence, $\bar{w} \equiv 0$ in this case.

Consider the second case when $\bar{w}_m(x) > 0$ for all $x \in \Omega$. For each $j \neq m$, consider the problem:

$$\begin{cases} -\Delta z(x) = \bar{\lambda} b_{jm} \bar{w}_m + \bar{\lambda} \sum_{k \neq m, j} b_{jk} \bar{w}_k + \bar{\lambda} b_{jj} z(x) \\ \quad + \bar{\lambda} k_j(\bar{w}_1, \dots, \bar{w}_{j-1}, z(x), \bar{w}_{j+1}, \dots, \bar{w}_n) z & \text{in } \Omega, \\ z = 0 & \text{on } \partial\Omega. \end{cases} \quad (4.1)$$

The function $z \equiv 0$ is a lower solution. From [A1], [A2] and $\partial \bar{w}_m / \partial \nu < 0$ on $\partial\Omega$, we have

$$\begin{aligned} \lambda_1 \delta \phi(x) &< \bar{\lambda} b_{jm}(x) \bar{w}_m + \bar{\lambda} \sum_{k \neq m, j} b_{jk} \bar{w}_k + \bar{\lambda} b_{jj} \delta \phi \\ &+ \bar{\lambda} k_j(\bar{w}_1, \dots, \bar{w}_{j-1}, \delta \phi, \bar{w}_{j+1}, \dots, \bar{w}_n) \delta \phi & \text{in } \Omega \end{aligned}$$

for $\delta > 0$ sufficiently small. Thus, $z = \delta \phi$ is a family of lower solution for problem (4.1) for such δ . Hence, the solution of (4.1) satisfies $\bar{w}_j(x) > \delta \phi(x)$ in Ω . This contradicts the fact that $\bar{w} \in \partial P$. We must thus have $\bar{w} \equiv 0$ in Ω in any case, for $\bar{\lambda} \geq 0$.

Next, define $z_i := w_i / \|w_i\|_E, i = 1, 2, \dots$; they satisfy

$$z_i + \lambda_i \Delta^{-1}[Bz_i] + \lambda_i \Delta^{-1}[G(w_i)w_i] / \|w_i\|_E = 0. \quad (4.2)$$

Since $\Delta^{-1}B$ is compact, there exists a subsequence (again denoted by $\{z_i\}$ for convenience) such that $\Delta^{-1}[Bz_i]$ converges in E . Since $\Delta^{-1}[G(w_i)w_i] / \|w_i\|_E$ tends to zero in E , as $\|w_i\|_E \rightarrow \|\bar{w}\|_E = 0$, equation (4.2) implies that $\{z_i\}$ converges in E to a function z_0 say, and

$$z_0 = -\bar{\lambda} \Delta^{-1}[Bz_0].$$

Moreover, we have $z_0 \in P$, since $w_i \in P$; and $z_0 \neq 0$ since $\|z_i\|_E = 1$. Hence, we must have $\bar{\lambda} > 0$. The uniqueness part Theorem 2.1 implies that $\bar{\lambda} = \lambda_0$. Consequently, we must have $(\bar{\lambda}, \bar{w}) = (\lambda_0, 0)$. This completes the proof of Theorem 4.1.

Let λ_1 be the principal eigenvalue of the Laplace operator on Ω with zero Dirichlet boundary data, and $\phi > 0$ be a positive eigenfunction.

Theorem 4.2. *Assume all the hypotheses in Theorem 4.1 and $b_{mm}(x) \equiv 0$ in assumption [A1]. Suppose further that*

[A3] *There exists constant vector $\vec{g} = \text{col}(g_1, \dots, g_n)$, $g_i > 0$, such that $B(x)\vec{g} > \lambda_1\vec{g}$ for all $x \in \Omega$.*

[A4] *For each $j \neq m$, the function k_j described in [A2] has the property $\limsup_{N \rightarrow \infty} N^{-1} \max\{k_j(u_1, \dots, u_{j-1}, N, u_{j+1}, \dots, u_n) : 0 \leq u_\tau \leq N, \tau \neq j\} < 0$. Then (1.1) has a steady-state solution in P for $\lambda = 1$.*

Proof. Let $\hat{z} = \phi\vec{g}$, [A3] implies that there exists $\mu > 1$ such that $B\hat{z} \geq \mu\lambda_1\hat{z}$ in $\bar{\Omega}$. Hence, $(-\Delta)^{-1}(B\hat{z}) \geq \mu\lambda_1(-\Delta)^{-1}(\phi\vec{g}) = \mu\phi\vec{g} = \mu\hat{z}$. By Theorem 2.5 in [5], we have $1/\hat{\lambda}_0 \geq \mu > 1$, that is $\hat{\lambda}_0 < 1$.

The equation for u_m and $b_{mm} \equiv 0$ imply that if $\lambda \in [0, C]$ for some C , and all the components $u_j, j \neq m$ of a steady state solution of (1.1) have the property $|u_j| \leq M$, then u_m must satisfy $|u_m| \leq KM$ for some constant K independent of M . We now show that there must indeed exist some constant M , such that if $\lambda \in (0, C]$ and $(\lambda, u) \in \mathcal{G}$, then $0 \leq u_j \leq M$ for all $j \neq m$. Otherwise, there exists $u_r, r \neq m$, and a point $x^* \in \Omega$, where $u_r(x^*) = \max_{j \neq m} \{\sup_{x \in \bar{\Omega}} u_j(x)\} > M$. Let $p = u_r(x^*)/M > 1$, then we have the property: $0 \leq u_j(x) \leq pM$ for all $x \in \bar{\Omega}, j \neq m$. Consider the equation satisfied by u_r at the point $x^* \in \Omega$.

$$\begin{aligned} -\Delta u_r(x^*) &= \sum_{j \neq m, r} \lambda b_{rj} u_j + \lambda b_{rm} u_m \\ &+ \lambda k_r(u_1, \dots, u_{r-1}(x^*), pM, u_{r+1}(x^*), \dots, u_n) u_r(x^*) \\ &\leq \lambda \left[\sum_{j \neq m, r} b_{rj} pM + b_{rm} K pM + \max\{(pM)^{-1} k_r(u_1, \dots, u_{r-1}, pM, u_{r+1}, \dots, u_n) \right. \\ &\quad \left. : 0 \leq u_\tau \leq pM, \tau \neq r\} (pM)^2 \right] < 0. \end{aligned}$$

The last inequality is satisfied for M sufficiently large due to hypothesis [A4]. This contradicts the definition of x^* . We thus assert that if $\lambda \in [0, C]$, then $|u_i|$ are bounded for each $i = 1, \dots, n$, if $(\lambda, u) \in \mathcal{G}$. Finally, using gradient estimates by means of (2.6), we conclude that \mathcal{G} cannot be unbounded if $\lambda \in [0, C]$, for some C . Thus, a solution of (2.6) must exist for $\lambda = 1$.

Example 4.1. We consider Example 3.2 with only a single modification. We change the entry b_{33} of B from -1 to 0. Clearly, condition [A1] is satisfied with $m = 3$. For [A2], we define

$$k_1(u_1, u_2) = \begin{cases} -u_1 - (u_2/u_1) \sin(u_1 u_2) & \text{if } u_1 \neq 0 \\ -u_2^2 & \text{if } u_1 = 0 \end{cases}$$

$$k_2(u_1, u_2) = \begin{cases} -(u_1/u_2) \sin(u_1 u_2) - u_2^2 & \text{if } u_2^2 \neq 0 \\ -u_1^2 & \text{if } u_2 = 0, \end{cases}$$

then k_1 and k_2 are continuous for $u_i \geq 0$. Condition [A3] is satisfied by choosing $\vec{g} = \text{col}(1,1,1)$. Using the formulas for k_1 and k_2 , we can verify [A4] is satisfied. Thus, Theorem 4.2 can be applied, and we conclude that (1.1) has a steady-state solution for $\lambda = 1$ in this case. As for Example 3.2 in the last section, we also find the bifurcating positive steady-state is asymptotically stable for λ close to $\hat{\lambda}_0$ for this example. The stability of the steady-state when $\lambda = 1$ remains to be considered.

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REFERENCES

- [1] M. Crandall and P. Rabinowitz, *Bifurcation, perturbation of simple eigenvalues and linearized stabilities*, Arch. Rat. Mech. Anal., 52 (1973), 161-181.
- [2] K. Deimling, "Nonlinear Functional Analysis," Springer, N.Y., 1985.
- [3] L. Evans, "Partial Differential Equations," AMS, Providence, 1998.
- [4] D. Henry, "Geometric Theory of Semilinear Parabolic Equations," Lecture Notes in Mathematics, Vol 840, Springer-Verlag, N.Y., 1981.
- [5] M. A. Kransnosel'skii, "Positive Solutions of Operator Equations," P. Noordhoff Ltd., Groningen, 1964.
- [6] A. Leung, "Systems of Nonlinear Partial Differential Equations: Applications to Biology and Engineering," MIA, Kluwer, Boston, 1989.
- [7] A. Leung and L.S. Ortega, *Bifurcating solutions and stabilities for multigroup neutron fission systems with temperature feedback*, J. Math. Anal. Appl., 194 (1995), 489-510.
- [8] A. Leung and B. Villa, *Reaction-diffusion systems for multigroup neutron fission with temperature feedback: positive steady-state and stability*, Diff. and Int. Eqs., 10 (1997), 739-756.
- [9] A. Pazy, "Semigroups of Linear Operators and Applications to Partial Differential Equations," Springer, N.Y., 1983.
- [10] M. C. Salvatori and E. Vitillaro, *Decay for the solutions of nonlinear abstract damped evolution equations with applications to partial and ordinary differential systems*, Diff. and Int. Eqs., 11 (1998), 223-262.
- [11] R. Temam, "Infinite Dimensional Dynamical Systems in Mechanics and Physics," Appl. Math. Sciences 68, Springer-Verlag, N.Y., 1997.
- [12] G. F. Webb, *A bifurcation problem for a nonlinear hyperbolic partial differential equation*, SIAM J. Math. Anal., 10 (1979), 922-932.