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On the functional solutions of a system of Partial Differential Equations relevant in mathematical physics

Abstract. We study the system of P.D.E.

$$abla \cdot \left[\sum_{j=1}^n a_{ij}(\boldsymbol{u},w) \nabla u_j + b_i(\boldsymbol{u},w) \nabla w\right] = 0, \ i = 1,\dots, n$$

$$\nabla \cdot (K(\boldsymbol{u}, w) \nabla w) = 0 \text{ in } \Omega$$

and the class of its solutions $(\boldsymbol{u}(\boldsymbol{x}), w(\boldsymbol{x})) = (u_1(\boldsymbol{x}), \dots, u_n(\boldsymbol{x}), w(\boldsymbol{x}))$ which occurs when a functional relation between $\boldsymbol{u}(\boldsymbol{x})$ and $w(\boldsymbol{x})$ of the form $\boldsymbol{u}(\boldsymbol{x}) = \boldsymbol{U}(w(\boldsymbol{x}))$ exists. If the solution satisfies constant boundary conditions $\boldsymbol{U}(w)$ is shown to exist and to satisfy a Volterra-Fredholm integral equation. An application to the thermistor problem and to a bifurcation problem of filtration in porous media is also given.

Keywords. Functional solutions, system of PDE, Thermistor problem, Volterra-Fredholm integral equation.

Mathematics Subject Classification (2000): 35B15, 35J66.

1 - Introduction

The autonomous boundary value problem (P)

$$(1.1) \quad \nabla \cdot \left[\sum_{j=1}^{n} a_{ij}(\boldsymbol{u}, w) \nabla u_j + b_i(\boldsymbol{u}, w) \nabla w \right] = 0 \text{ in } \Omega, \ \boldsymbol{u} = (u_1, \dots, u_n), \ i = 1, \dots, n$$

(1.2)
$$\nabla \cdot (K(\boldsymbol{u}, w) \nabla w) = 0 \text{ in } \Omega$$

Received: July 13, 2010; accepted in revised form: November 4, 2010.

(1.3)
$$\boldsymbol{u} = 0 \text{ on } \Gamma_0, \ \boldsymbol{u} = \bar{\boldsymbol{u}} \text{ on } \Gamma_1, \boldsymbol{q}_i \cdot \boldsymbol{n} = 0 \text{ on } \Gamma_2$$

(1.4)
$$w = 0 \text{ on } \Gamma_0, \ w = \bar{w} \text{ on } \Gamma_1, \ K(\boldsymbol{u}, w) \frac{\partial w}{\partial n} = 0 \text{ on } \Gamma_2,$$

where $\nabla w = (w_{x_1}, \dots, w_{x_N})$, Ω is an open and bounded subset of \mathbf{R}^N with a regular boundary Γ which consists of three disjoint hypersurfaces Γ_0 , Γ_1 and Γ_2 , $\bar{\mathbf{u}} = (\bar{u}_1, \dots, \bar{u}_n)$ is a constant vector of \mathbf{R}^n , the \mathbf{q}_i 's are given by

$$\boldsymbol{q}_i = \sum_{i=1}^n a_{ij}(\boldsymbol{u}, w) \nabla u_j + b_i(\boldsymbol{u}, w) \nabla w, \ i = 1, \dots, n$$

and n is the unit vector of the normal to Γ_2 , has many physical motivations. We recall, in particular, the equations of the electrical heating of a conductor, the so-called "thermistor problem" [9], [7]

(1.5)
$$\nabla \cdot (\kappa(u)\nabla u + \varphi \sigma(u)\nabla \varphi) = 0 \text{ in } \Omega$$

(1.6)
$$\nabla \cdot (\sigma(u)\nabla \varphi) = 0 \text{ in } \Omega$$

(1.7)
$$u = 0 \text{ on } \Gamma_0, \ u = \bar{u} \text{ on } \Gamma_1, \ \kappa(u) \frac{\partial u}{\partial n} = 0 \text{ on } \Gamma_2$$

(1.8)
$$\varphi = 0 \text{ on } \Gamma_0, \ \varphi = \bar{\varphi} \text{ on } \Gamma_1, \ \sigma(u) \frac{\partial \varphi}{\partial n} = 0 \text{ on } \Gamma_2,$$

where $u(\mathbf{x})$ and $\varphi(\mathbf{x})$, $\mathbf{x} = (x, y, z)$ are the temperature and the electric potential. The electric conductivity $\sigma(u)$ and the thermal conductivity $\kappa(u)$ are given functions of the temperature. Also the flow of an incompressible fluid in a porous medium, taking into account the Soret and Dufour effects, leads [1], [3] to a problem of the form (P). More precisely we have, with $\Gamma_2 = \emptyset$,

(1.9)
$$\nabla \cdot [\nabla u + S \nabla v + f(u, v) \nabla w] = 0 \text{ in } \Omega$$

(1.10)
$$\nabla \cdot [\nabla v + D\nabla u + q(u, v)\nabla w] = 0 \text{ in } \Omega$$

$$(1.11) \nabla \cdot (K\nabla w) = 0 \text{ in } \Omega$$

(1.12)
$$u = 0 \text{ on } \Gamma_0, \ u = \bar{u} \text{ on } \Gamma_1$$

$$(1.13) v = 0 ext{ on } \Gamma_0, v = \bar{v} ext{ on } \Gamma_1$$

(1.14)
$$w = 0 \text{ on } \Gamma_0, \ w = \bar{w} \text{ on } \Gamma_1.$$

Here u is the concentration, v the temperature and w the pressure. Equation (1.11) follows from the incompressibility assumption and the Darcy's law, (1.9) and (1.10) are balance equations, in which S and D are the Soret and Dufour coefficients.

If Ω is a subset of \mathbf{R}^2 , the equations (1.1) and (1.2) are invariant under all conformal mappings $\phi(z)$ such that $|\phi'(z)| \neq 0$, as noticed in the case of problem (1.5)-(1.7) in [6]. Suppose in problem (P)

$$\Gamma_2 = \emptyset$$
, $K = 1$, $n = 1$, $a_{11} = 1$, $b_1 = b(u)$, $\bar{w} = 1$

and consider the conformal transformation which maps Ω onto the annulus $\tilde{\Omega}$ of radii 1 and R>1. In $\tilde{\Omega}$ problem (1.2), (1.4) is immediately solved in polar coordinates ρ , θ by

$$(1.15) w = \frac{\log \rho}{\log R}.$$

Moreover, if we restrict our attention to radial solutions, equation (1.1) gives

(1.16)
$$\frac{d}{d\rho} \left(\rho \frac{du}{d\rho} \right) + b_u(u) \frac{du}{d\rho} \left(\frac{1}{\log R} \right) = 0.$$

From (1.15) we have $\rho = R^w$, thus (1.16) gives

(1.17)
$$\frac{d}{dw}\left(\frac{du}{dw}\right) + b_u(u)\frac{du}{dw} = 0,$$

where every reference to the domain Ω has disappeared and u is considered as a function of w. This elementary remark suggests the following

Definition 1. If $(\mathbf{u}(\mathbf{x}), w(\mathbf{x}))$ is a solution to problem (P) and there exist n functions $(\mathcal{U}_1(w), \ldots, \mathcal{U}_n(w))$ of class $C^1([0, \overline{w}])$ such that

$$u_1(\mathbf{x}) = \mathcal{U}_1(w(\mathbf{x})), \dots, u_n(\mathbf{x}) = \mathcal{U}_n(w(\mathbf{x}))$$

we say that $(\mathbf{u}(\mathbf{x}), w(\mathbf{x}))$ is a functional solution of problem (P).

In a previous work [3] the link between problem (P) and a two-point problem for a system of ordinary differential equations (TPP) is established. We prove in Section 2 that this correspondence between the functional solutions of (P) and the solutions of (TPP) is biunivocal. Moreover, problem (TPP) can be restated as a Volterra-Fredholm integral equation, for which a theorem of existence is proved in [3], when $a_{ij} = \delta_{ij}$. In Section 3 we extend that result to cover the case when $a_{ij}(\boldsymbol{u}, w)$ is an uniformly positive-definite matrix. Section 4 deals with an application to the thermistor problem and presents a slight generalization of the usual condition of existence and nonexistence [4]. Section 5 applies the method to a problem of bifurcation relevant in the flow of fluid in a porous medium using the Liapunoff-Schmidt technique for the corresponding two-point problem.

2 - The two-point problem associated with problem (P)

All functional solutions of problem (P) can be computed explicitly if we can find the solutions $(U(w), k), k = (k_1, \dots, k_n)$ of the following two-point problem (TPP)

(2.1)
$$\sum_{i=1}^{n} a_{ij}(\boldsymbol{U}, w) \frac{d\mathcal{U}_{i}}{dw} + b_{i}(\boldsymbol{U}, w) = k_{i}K(\boldsymbol{U}, w), \ i = 1, \dots, n, \ \boldsymbol{U} = (\mathcal{U}_{1}, \dots, \mathcal{U}_{n})$$

$$(2.2) U(0) = 0, \ U(\bar{w}) = \bar{\boldsymbol{u}}$$

and the solution of

$$(2.3) \Delta z = 0 \text{ in } \Omega$$

$$(2.4) z = 0 \text{ on } \Gamma_0, \ z = 1 \text{ on } \Gamma_1, \frac{\partial z}{\partial n} = 0 \text{ on } \Gamma_2$$

is known. First order ordinary differential equations depending on a parameter whose solutions must satisfy two boundary conditions have been considered, among others, in [14] and [13]. We assume

$$(2.5) K_1 \geq K(\boldsymbol{U}, w) \geq K_0 > 0, \text{ for all } \boldsymbol{U} \in \boldsymbol{R}^n, \text{ and for all } w \in [0, \bar{w}]$$

and denote with $\{TPP\}$ the set of all solutions of problem (2.1), (2.2) and with $\{P\}$ the set of all functional solutions of problem (1.1)-(1.4). We claim that

$$(2.6) card\{TPP\} = card\{P\}.$$

Let the map $\mathcal{I}: \{TPP\} \to \{P\}$ be defined as follows. Suppose $(U(w), \mathbf{k}) \in \{TPP\}$, and consider the nonlinear boundary value problem in the unknown $w(\mathbf{x})$

(2.7)
$$\nabla \cdot (K(U(w), w) \nabla w) = 0 \text{ in } \Omega$$

(2.8)
$$w = 0 \text{ on } \Gamma_0, \ w = \bar{w} \text{ in } \Gamma_1, \ \frac{\partial w}{\partial n} = 0 \text{ in } \Gamma_2.$$

(2.7), (2.8) has one and only one solution. For, let

(2.9)
$$v = F(w) = \int_{0}^{w} K(\boldsymbol{U}(t), t) dt.$$

By (2.5) F maps one-to-one $[0, \bar{w}]$ onto $[0, F(\bar{w})]$. Moreover, under (2.9) problem (2.7), (2.8) transforms into

$$\Delta v = 0$$
 in Ω

$$v=0 \text{ on } \Gamma_0, \ v=F(\bar{w}) \text{ on } \Gamma_1, \ \frac{\partial v}{\partial n}=0 \text{ on } \Gamma_2$$

and

$$v(\mathbf{x}) = F(\bar{w})z(\mathbf{x}),$$

with z(x) given by (2.3), (2.4). Therefore, the only solution of problem (2.7)-(2.8) is

$$(2.10) w(\mathbf{x}) = F^{-1}(F(\bar{w})z(\mathbf{x})).$$

We claim that

(2.11)
$$(u(x), w(x)) = (U(w(x)), w(x)) \in \{P\}.$$

For, by (2.1) and (2.7) we have

$$\nabla \cdot \left[\left(\sum_{i=1}^{n} a_{ij}(\boldsymbol{U}(w), w) \frac{d\mathcal{U}_{i}}{dw} + b_{i}(\boldsymbol{U}(w), w) \right) \nabla w \right] = k_{i} \nabla \cdot (K(\boldsymbol{U}(w), w)) = 0.$$

Moreover, (2.11) satisfies all the boundary conditions in problem (P). Hence \mathcal{I} is well defined. On the other hand \mathcal{I} is one-to-one. For, let $(U(w), \mathbf{k}) \in \{TPP\}$, $(\tilde{U}(w), \tilde{\mathbf{k}}) \in \{TPP\}$, and $(\mathbf{u}(\mathbf{x}), w(\mathbf{x})) = \mathcal{I}(U(w), \mathbf{k})$, $(\tilde{\mathbf{u}}(\mathbf{x}), \tilde{w}(\mathbf{x})) = \mathcal{I}(\tilde{U}(w), \tilde{\mathbf{k}})$. If

$$\mathcal{I}(\boldsymbol{U}(w), \boldsymbol{k}) = \mathcal{I}(\tilde{\boldsymbol{U}}(w), \tilde{\boldsymbol{k}})$$

we have $u(x) = \tilde{u}(x)$ and $w(x) = \tilde{w}(x)$ and also $U(w) = \tilde{U}(w)$ since w(x) takes all values between 0 and \bar{w} . In addition we have: $k = \tilde{k}$. We prove that \mathcal{I} is onto. Let $(u(x), w(x)) = (U(w(x)), w(x)) \in \{P\}$ and define

(2.12)
$$\theta_i = \int_0^w \sum_{j=1}^n a_{ij}(\boldsymbol{U}(t), t) \frac{d\mathcal{U}_i}{dw}(t) + b_i(\boldsymbol{U}(t), t) dt$$

(2.13)
$$\zeta = \int_{0}^{w} K(\mathbf{U}(t), t) dt$$

$$\Theta_i(\mathbf{x}) = \theta_i(w(\mathbf{x})), \ Z(\mathbf{x}) = \zeta(w(\mathbf{x})).$$

By (1.1) and (1.2) we have

$$\Delta\Theta_i = 0$$
 in Ω , $\Delta Z = 0$ in Ω

$$\Theta_i=0,\; Z=0 \; ext{on} \; \Gamma_0, \; \Theta_i=C_i, \; Z=C \; ext{on} \; \Gamma_1, \; rac{\partial \Theta_i}{\partial n}=0, \; rac{\partial Z}{\partial n}=0 \; ext{on} \; \Gamma_2,$$

where the C_i and C are constants and $C \neq 0$ by (2.5). If $z(\mathbf{x})$ is given by (2.3), (2.4), we have $\Theta_i(\mathbf{x}) = C_i z(\mathbf{x})$, $Z(\mathbf{x}) = C_i z(\mathbf{x})$ and setting

$$(2.14) k_i = \frac{C_i}{C}$$

we have

$$(2.15) \Theta_i = k_i Z.$$

By (2.12), (2.13) and (2.15) we obtain

(2.16)
$$\int_{0}^{w} \left[\sum_{j=1}^{n} a_{ij}(\boldsymbol{U}(t), t) \frac{d\mathcal{U}_{i}}{dw} + b_{i}(\boldsymbol{U}(t), t) \right] dt = k_{i} \int_{0}^{w} K(\boldsymbol{U}(t), t) dt.$$

Taking the derivative of (2.16) with respect to w we obtain (2.1). Hence, with k given by (2.14), we have

$$\mathcal{I}(\boldsymbol{U}(w), \boldsymbol{k}) = (\boldsymbol{u}(\boldsymbol{x}), w(\boldsymbol{x})).$$

Remark 1. As an immediate consequence of (2.6) we can say that problem (P) has one and only one solution if and only if the same is true for problem (TPP). Moreover, the search of the functional solutions of problem (P) is broken into two steps: (i) to find the solutions (U(w), k) of problem (TPP) and (ii) to solve the problem

$$\Delta z = 0 \text{ on } \Omega, \ z = 0 \text{ on } \Gamma_0, \ z = 1 \text{ on } \Gamma_1, \ \frac{\partial z}{\partial n} = 0 \text{ on } \Gamma_2.$$

The solutions of problem (P) are then given by $(\mathbf{u}(\mathbf{x}), w(\mathbf{x})) = (\mathbf{U}(w(\mathbf{x})), w(\mathbf{x}))$, where $w(\mathbf{x})$ is given by (2.10). Step (ii) contains, so-to-speak, the geometrical part of (P) and step (i) the nonlinear part.

3 - Existence of solutions for the two-point problem (TPP)

In this section we prove the existence of at least one solution for the two-point problem

(3.1)
$$A(\boldsymbol{U}, w) \frac{d\boldsymbol{U}}{dw} = K(\boldsymbol{U}, w)\boldsymbol{k} - \boldsymbol{b}(\boldsymbol{U}, w), \ \boldsymbol{k} = (k_1, \dots, k_n)^T, \ \boldsymbol{b} = (b_1, \dots, b_n)^T$$

$$(3.2) U(0) = 0$$

$$(3.3) U(\bar{w}) = \tilde{\boldsymbol{u}}, \ \bar{w} > 0,$$

where $A(U, w) = (a_{ij}(U, w))$ is a symmetric $n \times n$ matrix continuous function of $U = (\mathcal{U}_1, \dots, \mathcal{U}_n)$ and of $w \in [0, \overline{w}]$ which satisfies

(3.4)
$$a_M |\xi|^2 \ge \sum_{ij=1}^n a_{ij}(U, w) \xi_i \xi_j \ge a_m |\xi|^2, \ a_m > 0, \ \forall \ \xi \in \mathbf{R}^n$$

for all $U \in \mathbb{R}^n$ and for all $w \in [0, \bar{w}]$. Moreover, we assume $K(U, w) \in C^0(\mathbb{R}^n \times [0, \bar{w}])$, $b(U, w) \in C^0(\mathbb{R}^n \times [0, \bar{w}]; \mathbb{R}^n)$ and

$$(3.5) K_1 > K(U, w) > K_0 > 0$$

$$|\boldsymbol{b}(\boldsymbol{U}, w)| \leq \beta_0 + \beta_1 |\boldsymbol{U}|^a, \ 0 \leq a < 1, \ \beta_1 > 0.$$

We transform (3.1)-(3.3) in a Volterra-Fredholm integral equation. By (3.4) A(U, w) is invertible and the inverse $C(u, w) = A^{-1}(U, w)$ satisfies

(3.7)
$$c_{M}|\xi|^{2} \geq \sum_{ij=1}^{n} c_{ij}(U, w)\xi_{i}\xi_{j} \geq c_{m}|\xi|^{2}, \ \forall \xi \in \mathbf{R}^{n}.$$

Thus we can rewrite (3.1) as follows

(3.8)
$$\frac{d\mathbf{U}}{dw} = E(\mathbf{U}, w)\mathbf{k} - C(\mathbf{U}, w)\mathbf{b}(\mathbf{U}, w),$$

where E(U, w) = K(U, w)C(U, w). Integrating (3.8) we obtain

(3.9)
$$U(w) = \int_{0}^{w} E(U(t), t)dt k - \int_{0}^{w} C(U(t), t) b(U(t), t)dt.$$

Recalling (3.3) we have

(3.10)
$$\int_{0}^{\bar{w}} E(\boldsymbol{U}(t), t) dt \boldsymbol{k} = \bar{\boldsymbol{u}} + \int_{0}^{\bar{w}} C(\boldsymbol{U}(t), t) \boldsymbol{b}(\boldsymbol{U}(t), t) dt.$$

Moreover, by (3.5) and (3.7)

$$\left(\int\limits_{0}^{\bar{w}}E(\boldsymbol{U}(t),t)dt\right)^{-1}$$

exists. Solving (3.10) with respect to k and substituting in (3.9) we arrive at

(3.11)
$$\boldsymbol{U}(w) = \int_{0}^{w} E(\boldsymbol{U}(t), t) dt \left(\int_{0}^{\bar{w}} E(\boldsymbol{U}(t), t) dt \right)^{-1} \left[\bar{\boldsymbol{u}} + \int_{0}^{\bar{w}} C(\boldsymbol{U}(t), t) \boldsymbol{b}(\boldsymbol{U}(t), t) dt \right]$$
$$- \int_{0}^{\bar{w}} C(\boldsymbol{U}(t), t) \boldsymbol{b}(\boldsymbol{U}(t), t) dt.$$

To find "a priori" estimates for the solutions of (3.11) we need the following

Lemma 1. Let $G(t) = (g_{ij}(t))$ be a symmetric and continuous matrix satisfying

(3.12)
$$g_M |\xi|^2 \ge \sum_{ij=1}^n g_{ij}(U, w) \xi_i \xi_j \ge g_m |\xi|^2$$
, for all $\xi \in \mathbb{R}^n$, $g_m > 0$.

Then

$$\left\| \left(\int_{0}^{w} G(t)dt \right) \left(\int_{0}^{\bar{w}} G(t)dt \right)^{-1} \right\| \leq \frac{g_{M}}{g_{m}}.$$

Proof. Let $\lambda_1(t) \leq \lambda_2(t) \leq \ldots \leq \lambda_n(t)$ be the eigenvalues of G(t). By (3.12) and the mini-max property of the eigenvalues we have $g_m \leq \lambda_1(t)$, $\lambda_n(t) \leq g_M$ and

(3.13)
$$\lambda_n(t)|\xi|^2 \ge \sum_{ij=1}^n g_{ij}(t)\xi_i\xi_j \ge \lambda_1(t)|\xi|^2, \text{ for all } \xi = \in \mathbf{R}^n.$$

Integrating (3.13) we have

$$\int\limits_0^w \lambda_n(t)dt |\xi|^2 \geq \sum_{ij=1}^n \int\limits_0^w g_{ij}(t)dt \xi_i \xi_j \geq \int\limits_0^w \lambda_1(t)dt |\xi|^2, \ \text{for all} \ \xi \in \pmb{R}^n.$$

If $\Lambda_1(w) \leq \Lambda_2(w) \leq \ldots \leq \Lambda_n(w)$ are the eigenvalues of the matrix $\int\limits_0^w G(t)dt$ we have

(3.14)
$$\int_{0}^{w} \lambda_{1}(t)dt \leq \Lambda_{1}(t), \ \Lambda_{n}(t) \leq \int_{0}^{w} \lambda_{n}(t)dt.$$

Hence, by (3.14)

$$\begin{split} & \left\| \int\limits_0^w G(t)dt \Big(\int\limits_0^{\bar{w}} G(t)dt \Big)^{-1} \right\| \leq \left\| \int\limits_0^w G(t)dt \right\| \, \left\| \Big(\int\limits_0^{\bar{w}} G(t)dt \Big)^{-1} \right\| \\ & \leq \frac{\varLambda_n(w)}{\varLambda_1(\bar{w})} \leq \frac{\int\limits_0^w \lambda_n(t)dt}{\int\limits_0^{\bar{w}} \lambda_1(t)dt} \leq \frac{g_M w}{g_m \bar{w}} \leq \frac{g_M}{g_m}. \end{split}$$

Theorem 1. If (3.4), (3.5) and (3.6) hold then the problem (3.1)-(3.3) has at least one solution.

Proof. Let $\mathcal{B} = C^0([0, \bar{w}]; \mathbf{R}^n)$. Define the operator $\mathcal{A} : \mathcal{B} \to \mathcal{B}$, $\mathcal{A} = (\mathcal{A}_1, \dots, \mathcal{A}_n)$

(3.15)
$$\mathcal{A}(\boldsymbol{U})(w) = \int_{0}^{w} E(\boldsymbol{U}(t), t) dt \left(\int_{0}^{\bar{w}} E(\boldsymbol{U}(t), t) dt \right)^{-1} \left[\bar{\boldsymbol{u}} + \int_{0}^{\bar{w}} C(\boldsymbol{U}(t), t) \boldsymbol{b}(\boldsymbol{U}(t), t) dt \right] - \int_{0}^{w} C(\boldsymbol{U}(t), t) \boldsymbol{b}(\boldsymbol{U}(t), t) dt.$$

Let $0 \le \mu \le 1$. We claim that all solutions of the equation

$$(3.16) U = \mu \mathcal{A}(U), \ U \in \mathcal{B}$$

are bounded in \mathcal{B} . For, let U(w) be any solution of (3.16). By Lemma 1 and (3.7) we have

$$|\boldsymbol{U}(w)| \leq \frac{K_1}{K_0} \frac{c_M}{c_m} |\bar{\boldsymbol{u}}| + \left(1 + \frac{K_1}{K_0} \frac{c_M}{c_m}\right) \int_0^{\bar{w}} c_M \left(\beta_0 + \beta_1 |\boldsymbol{U}(t)|^a\right) dt.$$

Hence

(3.17)
$$|U(w)| \le C_1 \left[1 + \int_0^w |U(t)|^a dt \right].$$

Let $m \ge 1$. Then by Hölder inequality we obtain from (3.17)

$$|oldsymbol{U}(w)| \leq C_2 igg[1 + ar{w}^{rac{m-a}{m}} igg(\int\limits_{0}^{ar{w}} |oldsymbol{U}(t)|^m dt igg)^{rac{a}{m}} igg]$$

and integrating from 0 to \bar{w}

$$\left(\int_{0}^{w}\left|\boldsymbol{U}(t)\right|^{m}dt\right)^{\frac{1}{m}}\leq\left[1+\left(\int_{0}^{w}\left|\boldsymbol{U}(t)\right|^{m}dt\right)^{\frac{a}{m}}\right].$$

By (3.6) we have

$$\left(\int\limits_{0}^{\bar{w}}\left|\boldsymbol{U}(t)\right|^{m}dt\right)^{\frac{1}{m}}\leq C_{4},$$

where C_4 does not depend on m. Hence

$$||U||_{\mathcal{B}} \leq C_4.$$

The operator \mathcal{A} is also compact and continuous. To see that, let $\{U^{(k)}(w)\}$ be a sequence bounded in \mathcal{B} and $V^{(k)} = \mathcal{A}(U^{(k)})$. From (3.5), (3.6) and Lemma 1 we obtain

that $\{V^{(k)}\}$ is also bounded in \mathcal{B}

On the other hand, $V_i^{(k)}$ is absolutely continuous and we have from (3.15)

$$\mathbf{V}^{(k)'}(w) = E(\mathbf{U}(w), w) \left(\int_{0}^{\bar{w}} E(\mathbf{U}(t), t) dt \right)^{-1} \left[\bar{\mathbf{u}} + \int_{0}^{\bar{w}} C(\mathbf{U}(t), t) \mathbf{b}(\mathbf{U}(t), t) dt \right]$$

$$-C(\mathbf{U}(w), w) \mathbf{b}(\mathbf{U}(w), w), \text{ a.e.}$$
(3.20)

By (3.18) the right hand side of (3.20) is bounded in \mathcal{B} . Hence by the Arzela's theorem the operator \mathcal{A} is compact. Therefore by the Schauder fixed point theorem in Schaefer's form [10] \mathcal{A} has a fixed point which gives a solution to problem (TPP). \square

4 - Example 1

We apply the present method to the case, quoted in the introduction, of the electrical heating of a conductor. We have

(4.1)
$$\nabla \cdot (\kappa(u)\nabla u + \varphi \sigma(u)\nabla \varphi) = 0 \text{ in } \Omega$$

$$(4.2) \nabla \cdot (\sigma(u)\nabla \varphi) = 0 \text{ in } \Omega$$

(4.3)
$$u = 0 \text{ on } \Gamma_0, \ u = \bar{u} \text{ on } \Gamma_1, \frac{\partial u}{\partial n} = 0 \text{ on } \Gamma_2$$

The associated two-point problem is the following

(4.5)
$$\kappa(\mathcal{U})\frac{d\mathcal{U}}{d\varphi} + \varphi\sigma(\mathcal{U}) = k\sigma(\mathcal{U})$$

$$(4.6) \mathcal{U}(0) = 0$$

$$\mathcal{U}(\bar{\varphi}) = \bar{u}.$$

Let us assume $\bar{\varphi} > 0$, $\bar{u} > 0$, $\kappa(u) \in C^0(\mathbb{R}^1)$, $\sigma(u) \in C^0(\mathbb{R}^1)$ and

$$\kappa(u) > 0, \ \sigma(u) > 0, \ u \in [0, \infty).$$

Define for u > 0

$$F(u) = \int_{0}^{u} \frac{\kappa(t)}{\sigma(t)} dt, \ L = F(\bar{u}), \ M = \int_{0}^{\infty} \frac{\kappa(t)}{\sigma(t)} dt.$$

By separation of variables, from (4.5) we have

(4.8)
$$F(\mathcal{U}) = -\frac{\varphi^2}{2} + k\varphi.$$

Using (4.7), we obtain from (4.8)

$$F(\mathcal{U}) = -rac{arphi^2}{2} + \left(rac{ar{arphi}}{2} + rac{L}{ar{arphi}}
ight)arphi.$$

Thus, if a solution of (4.5)-(4.7) exists, it can only be given by

$$\mathcal{U} = F^{-1}igg(-rac{arphi^2}{2} + igg(rac{ar{arphi}}{2} + rac{L}{ar{arphi}}igg)arphiigg).$$

Examining the graphs of $F(\mathcal{U})$ and of the parabola $f(\varphi) = -\frac{\varphi^2}{2} + \left(\frac{\bar{\varphi}}{2} + \frac{L}{\bar{\varphi}}\right)\varphi$, recalling that $\varphi \in [0, \bar{\varphi}]$, we have

Lemma 2. If $M=\infty$ or $M<\infty$ and $\bar{\varphi}^2<2L$, then problem (4.5)-(4.7) has one and only one solution. Let $M<\infty$ and $\bar{\varphi}^2>2L$. If

$$\frac{1}{8}\bar{\varphi}^2 + \frac{L}{2} + \frac{L^2}{4\bar{\varphi}^2} < M$$

then problem (4.5)-(4.7) has one and only one solution. If

$$\frac{1}{8}\bar{\varphi}^2 + \frac{L}{2} + \frac{L^2}{4\bar{\varphi}^2} \ge M$$

problem (4.5)-(4.7) has no solution.

When problem (4.5)-(4.7) has a solution the corresponding functional solution of (4.1)-(4.4) is $(\varphi(\mathbf{x}), u(\mathbf{x}))$, where $\varphi(\mathbf{x})$ is computed from

$$(4.9) \hspace{1cm} \nabla \cdot \left[\sigma \left(F^{-1} \left(-\frac{\varphi^2}{2} + \left(\frac{\bar{\varphi}}{2} + \frac{L}{\bar{\varphi}} \right) \varphi \right) \right) \nabla \varphi \right] = 0 \text{ in } \Omega$$

(4.10)
$$\varphi = 0 \text{ on } \Gamma_0, \ \varphi = \bar{\varphi} \text{ on } \Gamma_1, \ \frac{\partial \varphi}{\partial n} = 0 \text{ on } \Gamma_2$$

and

$$u(\mathbf{x}) = F^{-1} \left(-\frac{\varphi^2(\mathbf{x})}{2} + \left(\frac{\bar{\varphi}}{2} + \frac{L}{\bar{\varphi}} \right) \varphi(\mathbf{x}) \right).$$

In addition, the solution of (4.9), (4.10) can be expressed in terms of z(x), the solution of (2.3), (2.4). A question naturally arises: is the functional solution of problem (4.1)-(4.4) (or more generally of problem (P)) the only solution? The answer is affirmative [3], at least for problem (4.1)-(4.4).

5 - Example 2

Setting $S=0, D=0, \bar{w}=\lambda, \bar{u}=0, \bar{v}=0$ in (1.9)-(1.13) we arrive at the following nonlinear eigenvalue problem (P_1)

$$\nabla \cdot [\nabla u + \lambda f(u, v) \nabla z] = 0 \text{ in } \Omega, \ u = 0 \text{ on } \Gamma$$

$$\nabla \cdot [\nabla v + \lambda g(u, v) \nabla z] = 0 \text{ in } \Omega, \ v = 0 \text{ on } \Gamma$$

$$\Delta z = 0 \text{ in } \Omega$$

(5.2)
$$z = 0 \text{ on } \Gamma_0, \ z = 1 \text{ on } \Gamma_1, \ \frac{\partial z}{\partial n} = 0 \text{ on } \Gamma_2.$$

By (5.1) there is no loss in generality in assuming

$$f(0,0) = 0, g(0,0) = 0.$$

 (P_1) has for every λ the solution

$$(u(\mathbf{x}), v(\mathbf{x}), z(\mathbf{x})) = (0, 0, z(\mathbf{x})),$$

where $z(\mathbf{x})$ is given by (5.1), (5.2). Do exist branches of non-trivial solutions bifurcating from (5.3)?

The two-point problem (TPP_1) associated with (P_1) is

(5.4)
$$\frac{d\mathcal{U}}{dz} + \lambda f(\mathcal{U}, \mathcal{V}) = \gamma, \ \mathcal{U}(0) = 0, \ \mathcal{U}(1) = 0$$

(5.5)
$$\frac{d\mathcal{V}}{dz} + \lambda g(\mathcal{U}, \mathcal{V}) = \mu, \ \mathcal{V}(0) = 0, \ \mathcal{V}(1) = 0$$

or, equivalently,

$$\frac{d^2\mathcal{U}}{dz^2} + \lambda \left[f_u(\mathcal{U}, \mathcal{V}) \frac{d\mathcal{U}}{dz} + f_v(\mathcal{U}, \mathcal{V}) \frac{d\mathcal{V}}{dz} \right] = 0, \ \mathcal{U}(0) = 0, \ \mathcal{U}(1) = 0$$

$$\frac{d^2\mathcal{V}}{dz^2} + \lambda \left[g_u(\mathcal{U}, \mathcal{V}) \frac{d\mathcal{U}}{dz} + g_v(\mathcal{U}, \mathcal{V}) \frac{d\mathcal{V}}{dz} \right] = 0, \ \ \mathcal{V}(0) = 0, \ \ \mathcal{V}(1) = 0.$$

Let

$$F: D(F) \subseteq L^2((0,1); \mathbb{R}^2) \to L^2((0,1); \mathbb{R}^2),$$

 $D(F) = \mathbb{R}^1 \times H_0^1(0,1) \times H^2(0,1), \ U = (\mathcal{U}, \mathcal{V})$

and

$$\boldsymbol{F}(\lambda, \boldsymbol{U}) = (\mathcal{U}'' + \lambda [f_u(\mathcal{U}, \mathcal{V})\mathcal{U}' + f_v(\mathcal{U}, \mathcal{V})\mathcal{V}'], \mathcal{V}'' + \lambda [g_u(\mathcal{U}, \mathcal{V})\mathcal{U}' + g_v(\mathcal{U}, \mathcal{V})\mathcal{V}']).$$

Lemma 3. The differential of $F(\lambda, U)$

$$F'_{U}(\lambda, 0)[H] = H'' + \lambda AH', H = (\mathcal{H}, \mathcal{K}),$$

where

$$A = egin{bmatrix} f_u(0,0) & f_v(0,0) \ g_u(0,0) & g_v(0,0) \end{bmatrix}$$

is invertible if and only if A has purely imaginary proper values $\pm ib$, $b \neq 0$ and

(5.6)
$$\lambda = \lambda_n, \ \lambda_n = \frac{2\pi n}{h}, \ n \in \mathbb{N}.$$

The eigenvalues λ_n are all double with the eigenfunctions

(5.7)
$$\mathbf{y}_{1}^{(n)}(z) = \left(\sin\frac{2\pi}{b}nz, 1 - \cos\frac{2\pi}{b}nz\right),$$

$$\mathbf{y}_{2}^{(n)}(z) = \left(1 - \cos\frac{2\pi}{b}nz, -\sin\frac{2\pi}{b}nz\right).$$

Proof. Let us consider the two-point problem

(5.8)
$$\frac{d\mathbf{H}}{dz} + \lambda A\mathbf{H} = \mathbf{k}, \ \mathbf{k} = (k_1, k_2)^T$$

$$\mathbf{H}(0) = 0$$

$$\mathbf{H}(1) = 0.$$

The solution of the Cauchy problem (5.8), (5.9) is given by

$$m{H}(z) = \Big(e^{-\lambda Az}\int\limits_{0}^{z}e^{\lambda At}dt\Big)m{k}.$$

Thus, in view of condition (5.10), nontrivial solutions of problem (5.8)-(5.10) occur if and only if

$$det \Big(\int\limits_{0}^{1}e^{\lambda Jt}dt\Big)=0,$$

where J is the Jordan form of A. Of the four possible, the only case which needs to be considered is when

$$J = \begin{bmatrix} a & -b \\ b & a \end{bmatrix}.$$

We have

$$det \Big(\int\limits_0^1 e^{\lambda Jt} dt\Big) = rac{1}{\lambda^2 (a^2+b^2)} \Big[e^{2\lambda a} - 2e^{\lambda a} \cos{(\lambda b)} + 1\Big].$$

This is different from zero if $a \neq 0$. When a = 0, we have the countable sequence of double eigenvalues (5.6) with the eigensolutions (5.7).

By Lemma 3 bifurcation for problem (5.4)-(5.5) can be present only when the matrix A has purely imaginary eigenvalues and $\lambda = \lambda_n$. Let $\mathbf{H} = (\mathcal{H}, \mathcal{K})$ and define the operator

$$\boldsymbol{B}:D(\boldsymbol{B})\subseteq L^2((0,1);\boldsymbol{R}^2)\to L^2((0,1);\boldsymbol{R}^2),\ \boldsymbol{B}(\boldsymbol{H})=\Big(\mathcal{H}''-\frac{2\pi}{b}n\mathcal{K}',\mathcal{K}''+\frac{2\pi}{b}n\mathcal{H}'\Big).$$

Lemma 4. B is symmetric.

Proof. Let $\tilde{\boldsymbol{H}} = (\tilde{\mathcal{H}}, \tilde{\mathcal{K}}) \in D(\boldsymbol{B})$. We have, integrating by parts,

$$\int_{0}^{1} \boldsymbol{B}(\boldsymbol{H}) \cdot \tilde{\boldsymbol{H}} dz = \int_{0}^{1} \mathcal{H} \left(\tilde{\mathcal{H}}'' - \frac{2\pi}{b} n \tilde{\mathcal{K}}' \right) + \mathcal{K} \left(\tilde{\mathcal{K}}'' + \frac{2\pi}{b} n \tilde{\mathcal{H}}' \right) dz = \int_{0}^{1} \boldsymbol{H} \cdot \boldsymbol{B}(\tilde{\boldsymbol{H}}) dz$$

(here the dot denotes the scalar product in \mathbb{R}^2).

Hereafter we examine the specific case of a quadratic nonlinear terms i.e. we assume

$$f(u,v) = -bv + A_1u^2 + A_2uv + A_3v^2$$
, $g(u,v) = bu + B_1u^2 + B_2uv + B_3v^2$.

The relevant two-point problem is

$$\mathcal{U}'' + \frac{2\pi}{b}n\left[-b\mathcal{V}' + A_1\mathcal{U}\mathcal{U}' + A_2(\mathcal{U}\mathcal{V}' + \mathcal{V}\mathcal{U}') + A_3\mathcal{V}\mathcal{V}'\right] = 0$$

$$\mathcal{V}'' + \frac{2\pi}{b}n\left[b\mathcal{U}' + B_1\mathcal{U}\mathcal{U}' + B_2(\mathcal{U}\mathcal{V}' + \mathcal{V}\mathcal{U}') + B_3\mathcal{V}\mathcal{V}'\right] = 0$$

$$\mathcal{U}(0) = \mathcal{V}(0) = 0, \ \mathcal{U}(1) = \mathcal{V}(1) = 0.$$

We apply the Liapunoff-Schmidt procedure and the Newton method [11], [12], [8], [5] to assess if bifurcation occurs or not. In view of the symmetry of B the nullspace of the adjoint B^* is spanned by $y_1^{(n)}$ and $y_2^{(n)}$. Let

$$G(U,\lambda) = \left(\mathcal{U}'' + \lambda [-b\mathcal{V}' + A_1\mathcal{U}\mathcal{U}' + A_2(\mathcal{U}\mathcal{V}' + \mathcal{V}\mathcal{U}') + A_3\mathcal{V}\mathcal{V}'],$$

$$\mathcal{V}'' + \lambda [b\mathcal{U}' + B_1\mathcal{U}\mathcal{U}' + B_2(\mathcal{U}\mathcal{V}' + \mathcal{V}\mathcal{U}') + B_3\mathcal{V}\mathcal{V}'] \right), \ U = (\mathcal{U}, \mathcal{V})$$

and

$$G_{10} = G'_{u}, \ G_{11} = \frac{1}{2!}G''_{u\lambda}, \ G_{20} = \frac{1}{2!}G''_{uu}, \ G_{21} = \frac{1}{3!}G'_{uu\lambda}.$$

Let $\mu = \lambda - \lambda_n$. The Liapunoff-Schmidt equations corresponding to problem (5.4) are

$$(5.11) \int_{0}^{1} \left(\boldsymbol{G}_{20}(0, \lambda_{n})[\boldsymbol{U}, \boldsymbol{U}] \cdot \boldsymbol{y}_{1}^{(n)} + \boldsymbol{G}_{11}(0, \lambda_{n})[\boldsymbol{U}, \mu] \cdot \boldsymbol{y}_{1}^{(n)} + \boldsymbol{G}_{21}(0)[\boldsymbol{U}, \boldsymbol{U}, \mu] \cdot \boldsymbol{y}_{1}^{(n)} \right) dz = 0$$

$$(5.12) \int_{0}^{1} \left(\boldsymbol{G}_{20}(\boldsymbol{0}, \lambda_{n}) [\boldsymbol{U}, \boldsymbol{U}] \cdot \boldsymbol{y}_{2}^{(n)} + \boldsymbol{G}_{11}(\boldsymbol{0}, \lambda_{n}) [\boldsymbol{U}, \mu] \cdot \boldsymbol{y}_{2}^{(n)} + \boldsymbol{G}_{21}(\boldsymbol{0}) [\boldsymbol{U}, \boldsymbol{U}, \mu] \cdot \boldsymbol{y}_{2}^{(n)} \right) dz = 0.$$

Set

$$(5.13) \hspace{1cm} \boldsymbol{U} = \xi_1 \varphi_1^{(n)} + \xi_2 \varphi_2^{(n)} + \text{higher order terms}.$$

Substituting (5.13) in (5.12) and in (5.11) we find

(5.14)
$$\Phi_{20}(\xi_1, \xi_2) + \Phi_{11}(\xi_1, \xi_2)\mu + o(|\xi|^i |\mu|^j) = 0, \ i+j > 2$$

(5.15)
$$\Psi_{20}(\xi_1, \xi_2) + \Psi_{11}(\xi_1, \xi_2)\mu + o(|\xi|^i |\mu|^j) = 0, \ i+j>2,$$

where

$$\begin{split} \varPhi_{20}(\xi_1,\xi_2) &= 2n\pi^2 \Big[\Big(A_3 - B_2 \Big) \xi_1^2 + \Big(A_1 + B_2 \Big) \xi_2^2 + \Big(B_3 - B_1 + 2A_2 \Big) \xi_1 \xi_2 \Big] \\ \varPhi_{11}(\xi_1,\xi_2) &= 2n\pi \xi_1 \end{split}$$

$$\begin{split} \varPsi_{20}(\xi_1,\xi_2) &= 2n\pi^2 \Big[\Big(-A_2 - B_2 \Big) \xi_1^2 + \Big(A_2 - B_1 \Big) \xi_2^2 + \Big(A_3 - 2B_2 - A_1 \Big) \xi_1 \xi_2 \Big] \\ \varPsi_{11}(\xi_1,\xi_2) &= -2n\pi\xi_2. \end{split}$$

We are interested in those solutions (ξ_1, ξ_2) of (5.14), (5.15) which tend to zero as $\mu \to 0$ i.e. in small solutions. Therefore we need only to consider the decreasing part of the Newton diagram, which in this case reduces to the segment joining the points (1,1) and (2,0). Hence the number of small solutions is given by the points of intersections different from (0,0) in the plane (ξ_1,ξ_2) of the two conics

(5.16)
$$\varPhi_{20}^{(n)}(\xi_1, \xi_2) + \varPhi_{11}^{(n)}(\xi_1, \xi_2)\mu = 0$$

(5.17)
$$\Psi_{20}^{(n)}(\xi_1, \xi_2) + \Psi_{11}^{(n)}(\xi_1, \xi_2)\mu = 0.$$

From a generic point of view, we have either one or three branches of small solutions

starting from each eigenvalue. Consider, as an example, the problem

$$\Delta u + 2\pi \nabla \cdot ((v + u^2 - uv + v^2)\nabla z) = 0 \text{ in } \Omega, \ u = 0 \text{ on } \Gamma$$

$$\Delta v + 2\pi \nabla \cdot ((-u + u^2 - v^2)\nabla z) = 0 \text{ in } \Omega, \ v = 0 \text{ on } \Gamma$$

$$\Delta z = 0 \text{ in } \Omega, \ z = 0 \text{ on } \Gamma_0, \ z = 1 \text{ on } \Gamma_1, \ \frac{\partial z}{\partial u} = 0 \text{ on } \Gamma_2.$$

The conics (5.16), (5.17) are, in this case, the hyperbolas

$$\begin{split} \xi_2^2 + \xi_1^2 - 4\xi_1 \xi_2 - \frac{\mu}{\pi} \xi_1 &= 0 \\ -2\xi_2^2 + 2\xi_1^2 - \frac{\mu}{\pi} \xi_2 &= 0. \end{split}$$

The only point of intersection different from (0,0), is $\left(-\frac{\mu}{3\pi}, -\frac{2\mu}{3\pi}\right)$, to which it corresponds the branch of small solutions

$$u(\mathbf{x}) = -\frac{\mu}{3\pi} \sin 2\pi z(\mathbf{x}) - \frac{2\mu}{3\pi} (1 - \cos 2\pi z(\mathbf{x})) + o(|\mu|),$$

$$v(\mathbf{x}) = -\frac{\mu}{3\pi} (1 - 2\cos z(\mathbf{x})) + \frac{2\mu}{3\pi} \sin 2\pi z(\mathbf{x}) + o(|\mu|).$$

However, exceptional cases in which bifurcation does not occur are possible as in the following problem

(5.18)
$$\Delta u + 2\pi \nabla \cdot ((v+u^2)\nabla z) = 0 \text{ in } \Omega, \ u = 0 \text{ on } \Gamma$$

(5.20)
$$\Delta z = 0 \text{ in } \Omega, \ z = 0 \text{ on } \Gamma_0, \ z = 1 \text{ on } \Gamma_1, \ \frac{\partial z}{\partial n} = 0 \text{ on } \Gamma_2.$$

In this case the conics (5.16), (5.17) are the hyperbolae

$$\xi_2^2 + \xi_1 \xi_2 - \frac{\mu}{\pi} \xi_1 = 0$$
$$\xi_1^2 + \xi_1 \xi_2 \frac{\mu}{\pi} \xi_2 = 0$$

with parallel asymptotes. Their only point of intersection is (0,0). Thus the trivial solution of problem (5.18)-(5.20) is isolated.

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