# Optimal Design of a Minimum Weight Thermal Diffuser with Constraint on the Output Thermal Power Flux

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Abstract. The object of this paper is the development of efficient mathematical and numerical tools to find the optimal shape of a minimum-weight thermal diffuser with a priori specifications on the inward thermal power flux (TPF) and a bound on the outward TPF. The present problem arises in connection with the use of high-power solid state devices in future communications satellites. In a space application the thermal power must ultimately be dissipated to the environment by using heatpipes and correspondingly large radiating areas. However, heatpipes can accept only a limited TPF from a source. Hence we have the requirement of a minimum-weight thermal diffuser with a uniform bound on the outward TPF. Shape optimal design and finite elements methods are used. Complete numerical results are provided.

# 1. Introduction

The object of this paper is the development of efficient mathematical and numerical tools to find the optimal shape of a minimum-weight thermal diffuser with a priori specifications on the input and output thermal power flux. This paper contains a working theory, the necessary computations and numerical examples. A similar problem was studied by Ph. Destuynder [5] with the require-

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ment that the temperature at every point of the diffuser be less than a specified critical temperature.

The present problem arises in connection with the use of high-power solid state devices (HPSSD's) in future communication satellites. The specifications for this diffuser came from the Center for Research in Communication (CRC) in Canada<sup>1</sup>.

"An HPSSD dissipates a large amount of thermal power (typ. > 50W) over a relatively small mounting surface (typ.  $1.25 \text{cm}^2$ ). Yet its junction temperature is required to be kept moderately low (typ.  $110^{\circ}$ C). The thermal resistance from the junction to the mounting surface is known for any particular HPSSD (typ.  $1^{\circ}$ C/W), so that the mounting surface is required to be kept at a lower temperature than the junction (typ.  $60^{\circ}$ C).

In a space application the thermal power must ultimately be dissipated to the environment by the mechanism of radiation. However, to radiate large amounts of thermal power at moderately low temperatures, correspondingly large radiating areas are required. Thus we have the requirement to efficiently spread the high thermal power flux (TPF) at the HPSSD source (typ.  $40W/cm^2$ ) to a low TPF at the radiator (typ.  $.04W/cm^2$ ) so that the source temperature is maintained at an acceptably low level (typ.  $< 60^{\circ}$ C at mounting surface). The efficient spreading task is best accomplished using heatpipes, but the snag in the scheme is that heatpipes can accept only a limited maximum TPF from a source (typ.  $4W/cm^2$ ).

Hence we are led to the requirement for a thermal diffuser. This device is inserted between the HPSSD and the heatpipes, and reduces the TPF at the source (typ. >  $40W/cm^2$ ) to a level acceptable to the heatpipes (typ. max.  $4W/cm^2$ ). The heatpipes then sufficiently spread the heat over large space radiators, reducing the TPF from a level at the diffuser (typ. max.  $4W/cm^2$ ) to that at the radiator (typ.  $.04W/cm^2$ ). This scheme of heat spreading is depicted in Fig. 1.

It is the design of the thermal diffuser which is the problem at hand. We may assume that the HPSSD presents a uniform thermal power flux to the diffuser at the HPSSD/diffuser interface. Heatpipes are essentially isothermalizing devices, and we may assume that the diffuser/heatpipes interface is indeed isothermal. Any other surfaces of the diffuser may be treated as adiabatic."

Some early results were presented in Delfour-Payre-Zolésio [4].

**Notation.**  $\mathbb{R}$  (resp.  $\mathbb{R}_+$ ) is the field (resp. semigroup) of all (resp. positive or zero) reals. Given an integer  $n \ge 1$ ,  $\mathbb{R}^n$  is the *n*-dimensional Euclidean space. The topological dual of a Banach space E will be denoted by E'. Let  $\Omega$  be an open subset of  $\mathbb{R}^n$ .  $\mathfrak{N}(\Omega)$  is the vector space of all real functions defined on  $\Omega$  whose partial derivatives of all orders exist and are continuous and whose support is contained in some compact subset of  $\Omega$ . For  $k \ge 1$  an integer,  $H^k(\Omega)$  denotes the vector space of all real functions f on  $\Omega$  such that f and its distributional derivatives  $D^s f$  of order  $|s| = \sum_{i=1}^n |s_i| \le k$  all belong to  $L^2(\Omega)$ . For real  $k \ge 0$ ,

<sup>&</sup>lt;sup>1</sup>The statement of the problem and Figure 1 have been graciously provided by Dr. V. A. Wehrle of CRC.



Fig. 1. Heat-spreading scheme for high-power solid state devices.

 $H^{k}(\Omega)$  is defined by interpolation; for negative reals s < 0,  $H^{s}(\Omega) = (H_{0}^{-s}(\Omega))'$ where  $H_{0}^{s}(\Omega)$  is the closure of  $\mathfrak{N}(\Omega)$  in the  $H^{s}(\Omega)$ -topology.

Given a Banach space E,  $C^0(\Omega; E)$  is the vector space of continuous functions defined on  $\Omega$  into E. For  $k \ge 1$  an integer,  $C^k(\Omega; E)$  is the vector space of all continuous functions defined on  $\Omega$  into E whose partial derivatives are continuous up to order k. When  $E = \mathbb{R}$  we shall write  $C^0(\Omega)$  and  $C^k(\Omega)$ .

### 2. Statement of the Problem

We assume that the thermal diffuser is a volume  $\Omega$  symmetrical about the z-axis (cf. Figure 2A) whose boundary surface is made up of three regular pieces: the mounting surface  $\Sigma_1$  (a disk perpendicular to the z-axis with center in (r, z) = (0, 0)), the lateral adiabatic surface  $\Sigma_2$  and the interface  $\Sigma_3$  between the diffuser and the heatpipes saddle (a disk perpendicular to the z-axis with center in (r, z) = (0, L)).

The temperature distribution over this volume  $\Omega$  is the solution of the stationary heat equation  $\Delta T = 0$  ( $\Delta T$ , the Laplacian of T) with the following boundary conditions on the surface  $\Sigma = \Sigma_1 \cup \Sigma_2 \cup \Sigma_3$  (or boundary of  $\Omega$ ):

$$k \frac{\partial T}{\partial n} = q_{in} \quad \text{on } \Sigma_1$$

$$k \frac{\partial T}{\partial n} = 0 \quad \text{on } \Sigma_2$$

$$T = T_3 \quad \text{on } \Sigma_3, T_3 = \text{constant},$$

$$(2.1)$$

where *n* always denotes the outward normal to the boundary surface  $\Sigma$  and

 $\partial T/\partial n$  is the normal derivative to the boundary surface  $\Sigma$ ,

$$\frac{\partial T}{\partial n} = \nabla T \cdot n \qquad (\nabla T = \text{gradient of } T). \tag{2.2}$$

The parameters appearing in (2.1) are:

- k = thermal conductivity (typ. 1.8W/cm.×°C)
- $q_{\rm in}$  = uniform inward thermal power flux at the source (positive constant).

The radius  $R_0$  of the mounting surface  $\Sigma_1$  is fixed so that the boundary surface  $\Sigma_1$  is already given in the design problem.

For practical considerations, we assume the diffuser to be solid without interior hollows or cutouts. The class of shapes for the diffuser is characterized by the design parameter  $L \ge 0$  and the positive function R(z),  $0 \le z \le L$ , with  $R(0) = R_0 > 0$ . They are volumes of revolution  $\Omega$  about the z-axis generated by the surface A between the z-axis and the function R(z) (cf. Fig. 2B), that is

$$\Omega = \{ (x, y, z) | 0 \le z \le L, x^2 + y^2 \le R^2(z) \}.$$
(2.3)

So the shape of  $\Omega$  is completely specified by the length L and the function R on the interval [0, L].

Assuming that the diffuser is made up of a homogeneous material of uniform density (no hollows) the design objective is to minimize the volume

$$J(\Omega) = \pi \int_0^L R^2(z) \, dz$$
 (2.4)

subject to a uniform constraint on the outward thermal power flux at the interface  $\Sigma_3$  between the diffuser and the heatpipes saddle:

$$\sup_{p \in \Sigma_3} -k \frac{\partial T}{\partial z}(p) \leqslant q_{\text{out}},$$
(2.5)

where  $q_{out}$  is a specified positive constant.



Fig. 2. (A) Volume  $\Omega$  and its boundary  $\Sigma$ ; (B) Surface A generating  $\Omega$ ; (C) Surface D generating  $\overline{\Omega}$ .

It is readily seen that the minimization problem (2.4) subject to the constraint (2.5) (where T is the solution of the heat equation with the boundary conditions (2.1)) is independent of the fixed temperature on boundary  $\Sigma_3$ . In other words the optimal shape  $\Omega$ , if it exists is independent of  $T_3$ . As a result, from now on we set  $T_3$  equal to 0.

### 3. Reformulation of the Problem

In this problem, the shape parameter L and the shape function R are not independent of each other since the function R is defined on the interval [0, L]. This motivates the following change of variable on the z-axis

$$\zeta = \frac{z}{L}, \qquad 0 \le \zeta \le 1. \tag{3.1}$$

The length of the diffuser is now one and it is possible to work on a fixed interval [0, 1]. Similarly we can scale the  $x_1$  and  $x_2$  variables by  $R_0$ :

$$\xi_1 = \frac{x_1}{R_0}, \qquad \xi_2 = \frac{x_2}{R_0}. \tag{3.2}$$

We shall see that the shape parameter  $\tilde{L} = L/R_0$  will become a design parameter in the differential equation. The shape of the transformed domain is now completely specified by the new *shape function*  $\tilde{\rho}$  defined on the interval [0, 1] with the conditions  $\tilde{\rho}(0) = 1$  and  $\tilde{\rho}(\zeta) \ge 0$ ,  $0 \le \zeta \le 1$ . The former and new shape functions are related as follows:

$$\tilde{\rho}(\zeta) = R(L\zeta)/R_0. \tag{3.3}$$

Denote by  $\tilde{\Omega}$  and  $\tilde{\Sigma}$ ,  $\tilde{\Sigma}_1$ ,  $\tilde{\Sigma}_2$ ,  $\tilde{\Sigma}_3$  the transformed domain  $\Omega$  and surfaces  $\Sigma$ ,  $\Sigma_1$ ,  $\Sigma_2$ ,  $\Sigma_3$  in the new coordinate system ( $\xi_1$ ,  $\xi_2$ ,  $\zeta$ ) through the transformations (3.1) and (3.2). The *transformed domain*  $\tilde{\Omega}$  is solely dependent on the function  $\tilde{\rho}$  and is generated by the revolution of the surface D,

$$D = \langle (\rho, \zeta) | 0 < \zeta < 1, 0 < \rho < \tilde{\rho}(\zeta) \rangle, \qquad (3.4)$$

about the  $\zeta$ -axis (see Figure 2C). The surface D has four regular boundaries:

$$S_{1} = \langle (\rho, \zeta) | 0 < \rho < 1, \zeta = 0 \rangle$$

$$S_{2} = \langle (\rho, \zeta) | 0 < \zeta < 1, \rho = \tilde{\rho}(\zeta) \rangle$$

$$S_{3} = \langle (\rho, \zeta) | 0 < \rho < \tilde{\rho}(1), \zeta = 1 \rangle$$

$$S_{4} = \langle (\rho, \zeta) | 0 < \zeta < 1, \rho = 0 \rangle.$$
(3.5)

#### 3.1. Equations for the Scaled Temperature

Introduce the scaled temperature

$$y(\xi_1, \xi_2, \zeta) = \frac{k}{Lq_{\rm in}} T(R_0 \xi_1, R_0 \xi_2, L\zeta)$$
(3.6)

in the new variables  $(\xi_1, \xi_2, \zeta)$  or in cylindrical coordinates

$$y(\rho,\zeta) = \frac{k}{Lq_{\rm in}} T(R_0\rho, L\zeta), \qquad \rho = \sqrt{\xi_1^2 + \xi_2^2}.$$
 (3.6a)

This scaling of the original temperature is motivated by the fact that we want to define y in such a way that

$$\frac{k}{q_{\rm in}}\frac{\partial T}{\partial z}(r,L) = \frac{\partial y}{\partial \zeta}(\rho,1), \qquad 0 \le r = R_0 \rho \le R(L) = R_0 \tilde{\rho}(1). \tag{3.7}$$

This quantity appears in the constraint on the outward thermal power flux through the boundary surface  $\Sigma_3$ .

The solution of the heat equation in  $\Omega$  subject to boundary conditions (2.1) with  $T_3 = 0$  coincides with the solution of the following variational equation: to find T in  $H_0(\Omega)$  such that

$$k \int_{\Omega} \left[ \frac{\partial T}{\partial x_1} \frac{\partial v}{\partial x_1} + \frac{\partial T}{\partial x_2} \frac{\partial v}{\partial x_2} + \frac{\partial T}{\partial z} \right] d\Omega = q_{\rm in} \int_{\Sigma_1} v \, d\Sigma_1$$
(3.8)

for all v in  $H_0(\Omega)$ , where

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$$H_0(\Omega) = \left\{ v \in H^1(\Omega) : v|_{\Sigma_3} = 0 \right\}$$
(3.9)

is a closed linear subspace of  $H^1(\Omega)$ .

In the new variables  $(\xi_1, \xi_2, \zeta)$  the variational problem (3.8)–(3.9) becomes: to find y in  $H_0(\Omega)$  such that

$$\int_{\tilde{\Omega}} \left[ \left( \frac{L}{R_0} \right)^2 \left( \frac{\partial y}{\partial \xi_1} \frac{\partial v}{\partial \xi_1} + \frac{\partial y}{\partial \xi_2} \frac{\partial v}{\partial \xi_2} \right) + \frac{\partial y}{\partial \zeta} \frac{\partial v}{\partial \zeta} \right] d\tilde{\Omega} = \int_{\tilde{\Sigma}_1} v \, d\tilde{\Sigma}_1 \tag{3.10}$$

for all v in  $H_0(\bar{\Omega})$  (same definition as in (3.9)). Equivalently, it is easy to show that y is the solution of the following boundary value problem:

$$A(y) = -\left[\left(\frac{L}{R_0}\right)^2 \left(\frac{\partial^2 y}{\partial \xi_1^2} + \frac{\partial^2}{\partial \xi_2^2}\right) + \frac{\partial^2 y}{\partial \zeta^2}\right] = 0 \text{ in } \tilde{\Omega}$$
(3.11)

with boundary conditions

$$\frac{\partial y}{\partial \nu_A}\Big|_{\tilde{\Sigma}_1} = 1, \qquad \frac{\partial y}{\partial \nu_A}\Big|_{\tilde{\Sigma}_2} = 0, \qquad y|_{\tilde{\Sigma}_3} = 0, \qquad (3.12)$$

where

$$\frac{\partial y}{\partial \nu_{A}} = \left(\frac{L}{R_{0}}\right)^{2} \left(\nu_{1} \frac{\partial y}{\partial \xi_{1}} + \nu_{2} \frac{\partial y}{\partial \xi_{2}}\right) + \nu_{\zeta} \frac{\partial y}{\partial \zeta}$$
(3.13)

is the conormal derivative on the boundary  $\tilde{\Sigma}$  associated with the operator A and the unit outward normal vector  $v = (v_1, v_2, v_{\zeta})$  in the  $(\xi_1, \xi_2, \zeta)$  coordinate system. The cost function becomes

$$J(\Omega) = R_0^3 J(\tilde{L}, \tilde{\rho}), \qquad J(\tilde{L}, \tilde{\rho}) = \tilde{L}\pi \int_0^1 \tilde{\rho}(\zeta)^2 d\zeta, \qquad \tilde{L} = L/R_0, \qquad (3.14)$$

where  $\tilde{L}$  is now a dimensionless length; the constraint (2.5) reduces to

$$\sup_{p \in \tilde{\Sigma}_{3}} -\frac{\partial y}{\partial \nu_{\mathcal{A}}}(p) = \sup_{p \in \tilde{\Sigma}_{3}} -\frac{\partial y}{\partial \zeta}(p) \leq \frac{q_{\text{out}}}{q_{\text{in}}},$$
(3.15)

which only depends on the dimensionless flux ratio

$$q = q_{\rm out}/q_{\rm in}.\tag{3.16}$$

### 3.2. Variational Equations in Cylindrical Coordinates

In cylindrical coordinates the variational problem (3.8) becomes: to find T in  $V_0(A)$  such that

$$k \int_{0}^{L} \int_{0}^{R(z)} \left[ \frac{\partial T}{\partial r} \frac{\partial v}{\partial r} + \frac{\partial T}{\partial z} \frac{\partial v}{\partial z} \right] r \, dr \, dz = q_{\rm in} \int_{0}^{R_0} v r \, dr \tag{3.17}$$

for all v in  $V_0(A)$ , where

$$V_0(A) = \left\{ v | \sqrt{r} v, \sqrt{r} \frac{\partial v}{\partial r}, \sqrt{r} \frac{\partial v}{\partial z} \in L^2(A) \text{ and } v(r, L) = 0, 0 \leq r \leq R(L) \right\}.$$
(3.18)

Similarly the variational problem (3.10) for the scaled temperature y becomes: to find y in  $V_0(D)$  such that

$$\int_{0}^{1} \int_{0}^{\hat{\rho}(\zeta)} \left[ \tilde{L}^{2} \frac{\partial y}{\partial \rho} \frac{\partial v}{\partial \rho} + \frac{\partial y}{\partial \zeta} \frac{\partial v}{\partial \zeta} \right] \rho \, d\rho \, d\zeta = \int_{0}^{1} v(\rho, 0) \rho \, d\rho \tag{3.19}$$

for all v in  $V_0(D)$  (same definition as in (3.18), but in new coordiantes  $(\rho, \zeta)$ ).

#### 3.3. The Constrained Minimization Problem (P)

Denote by (P) the constrained minimization problem which consists in minimizing  $J(\tilde{L}, \tilde{\rho})$  in (3.14) with respect to the dimensionless design parameter  $\tilde{L}$  and design function  $\tilde{\rho}$  subject to the constraint

$$\sup\left\{-\frac{\partial y}{\partial \zeta}(\rho,1): 0 \le \rho \le \tilde{\rho}(1)\right\} \le q, \qquad (3.20)$$

where y is the solution of the variational problem (3.19).

It is readily seen that the optimal design is only a function of  $q = q_{out}/q_{in}$ . The parameter  $R_0$  only appears as a scaling parameter. If  $\tilde{L}^*$  and  $\tilde{\rho}^*$  are the optimal scaled parameter and function, then the optimal parameter  $L^*$  and function  $R^*$  for the original problem (2.1)-(2.5) are given by

$$L^* = R_0 \tilde{L}^*, \qquad R^*(z) = R_0 \tilde{\rho}^* (z/R_0 \tilde{L}^*).$$
(3.21a)

This corresponds to the following coordinate transformation:

$$x_1 = R_0 \xi_1, \qquad x_2 = R_0 \xi_2, \qquad z = R_0 \tilde{L}^* \zeta = L^* \zeta.$$
 (3.21b)

### 3.4. Angles and Smoothness of the Solution

In the forthcoming sections of this paper we shall need sufficient smoothness of the solution T to the heat equation with boundary conditions (2.1). Since the boundary conditions on each piece of boundary is constant, the global smoothness of the solution will only be affected by the smoothness of the function R(z)and the angles  $\theta_1$  and  $\theta_3$  between surfaces  $\Sigma_1$  and  $\Sigma_2$  and  $\Sigma_3$  and  $\Sigma_2$ , respectively (cf. Fig. 3)

In this section we assume that R (resp.  $\rho$ ) is sufficiently regular, so that the smoothness of T is solely affected by the angles  $\theta_1$  and  $\theta_3$  between the curve (z, R(z)) and the planes z = 0 and z = L, respectively.

Away from the curves  $C_1 = \Sigma_1 \cap \Sigma_2$  and  $C_3 = \Sigma_2 \cap \Sigma_3$ , the solution T is infinitely continuously differentiable. So we only need to study the smoothness of T in neighborhoods  $V_1$  and  $V_3$  of  $C_1$  and  $C_3$ , respectively.

In the special geometry where  $\theta_1 = \pi/2$  (resp.  $\theta_3 = \pi/2$ ), we can use the "principle of symmetry" to show that the solution T belongs to  $C^{\infty}(V_1)$  (resp.  $C^{\infty}(V_3)$ ) in a neighborhood  $V_1$  (resp.  $V_3$ ) of the curve  $C_1$  (resp.  $C_3$ ).



Fig. 3. Angles  $\theta_1$  and  $\theta_3$  between the surface  $\Sigma_2$  and the surfaces  $\Sigma_1$  and  $\Sigma_3$ .



Fig. 4. Smoothness of the solution in the vicinity of a corner as a function of the angle.

When the angles  $\theta_1$  and  $\theta_3$  are between 0 and  $\pi$  we can use the following results from P. Grisvard [7, 8, 9, 10]. We consider the two cases of Fig. 4.

In a neighborhood  $V_1$  of the curve  $C_1$ , the solution belongs to  $H^{1+(\pi/\theta_1)-\epsilon}(V_1)$ for all  $\epsilon > 0$ ; in a neighborhood  $V_3$  of the curve  $C_3$ , the solution belongs to  $H^{1+(\pi/2\theta_3)-\epsilon}(V_3)$  for all  $\epsilon > 0$ . Therefore, for the range of angles we are interested in, there exist sufficiently small  $\epsilon > 0$  and a neighborhood  $V_1$  (resp.  $V_3$ ) such that

$$\forall \theta_1, 0 < \theta_1 < \pi, \exists \varepsilon > 0, 1 + \frac{\pi}{\theta_1} - \varepsilon > 2$$

(resp.  $\forall \theta_3, 0 < \theta_3 < \pi, \exists \varepsilon > 0, 1 + \pi/2\theta_3 - \varepsilon > 3/2$ ). Then

$$T|_{V_1} \in H^{2+\sigma}(V_1) \quad \text{for } 0 \leq \sigma < \varepsilon$$
 (3.22)

and

$$T|_{V_2} \in H^{3/2+\sigma}(V_2) \quad \text{for } 0 \leq \sigma < \varepsilon.$$
(3.23)

### 4. Approximation of the Solution to the Constrained Minimization Problem (P)

In the absence of existence and uniqueness result, we shall assume the existence of at least one solution to problem (P) and concentrate on the approximation of the solution.

The constraint (3.20) is completely equivalent to the new constraint

$$f(\tilde{L},\tilde{\rho}) = 0, \tag{4.1}$$

where

$$f(\tilde{L},\tilde{\rho}) = \int_{\tilde{\Sigma}_3} \left[ \frac{\partial y}{\partial \zeta} + q \right]^- d\sigma = 2\pi \int_0^{\tilde{\rho}(1)} \left[ \frac{dy}{\partial \zeta}(\rho,1) + q \right]^- \rho \, d\rho \tag{4.2}$$

with  $u^{-} = \sup\{-u, 0\}$ .

Associate with an arbitrary family  $\{\varepsilon : \varepsilon > 0\}$  of small positive numbers the penalized cost function

$$J_{\varepsilon}(\tilde{L},\tilde{\rho}) = J(\tilde{L},\tilde{\rho}) + \frac{1}{\varepsilon}f(\tilde{L},\tilde{\rho}).$$
(4.3)

Replace the original constrained minimization problem (P) specified in section

3.3 by the following family of  $\epsilon$ -indexed unconstrained minimization problems  $(P_{\epsilon})_{\epsilon>0}$ : to find  $(\tilde{L}_{\epsilon}, \tilde{\rho}_{\epsilon})$  such that

$$J_{\varepsilon}(\tilde{L}_{\varepsilon}, \tilde{\rho}_{\varepsilon}) \leqslant J_{\varepsilon}(\tilde{L}, \tilde{\rho}) \tag{4.4}$$

for all  $\tilde{L}$  and  $\tilde{\rho}$  such that

$$\tilde{L} \ge 0$$
 and  $\forall \zeta$ ,  $\tilde{\rho}(\zeta) \ge 0$  with  $\tilde{\rho}(0) = 1$ . (4.5)

It is readily seen that any limit point  $\tilde{L}$ ,  $\tilde{\rho}$  of a sequence  $(\tilde{L}_{\epsilon_n}, \tilde{\rho}_{\epsilon_n})$  as  $\epsilon_n$  goes to zero is a global minimum solution to problem (P). This is a consequence of the fact that the function f is nonnegative.

### 5. Differentiability of the Scaled Temperature with Respect to the Shape

## 5.1. The Speed Method

We make use of the techniques introduced by J. Cea [1, 2] (cf. J.-P. Zolésio [12–15]). Let V be a time-dependent regular vector field defined in a neighborhood U of  $\overline{D}$  ( $\overline{D}$ , the closure of the domain D in Figure 2C) and for any point x in D consider the solution x to the differential equation

$$\frac{d}{dt}x(t,X) = V(t,x(t,X)), \qquad x(0,X) = X.$$
(5.1)

Then there exists a number  $t_1 > 0$  such that, for all X in D, the solution x(t, X) is defined for all t in  $[0, t_1]$ . Consider the transformation

 $T_t(V): X \to x(t, X).$ 

It changes the domain D into a new domain

$$D_t = T_t(V)(D), \quad 0 \le t < t_1,$$
 (5.2)

where

$$D_t = \{ x | x = x(t, X), X \in D \}.$$

It can be shown that if V is chosen in  $C^0(0,1; C^k(U; \mathbb{R}^3))$ , the transformation  $T_t(V)$  is a  $C^k$ -diffeomorphism from a neighborhood of D into its ring (cf. J.-P. Zolésio [12, 14]). The field is called the *deformation speed*.

### 5.2. Choice of the Speed V

The initial three-dimensional problem consists in specifying the domain  $\Omega$ . But  $\Omega$  is z-axisymmetrical and the scaled domain

$$\tilde{\Omega} = \left\{ (\xi_1, \xi_2, \zeta) | \xi_1 = \frac{x_1}{R_0}, \xi_2 = \frac{x_2}{R_0}, \zeta = \frac{z}{L}, (x_1, x_2, z) \in \Omega \right\}$$
(5.3)

is completely specified by its boundary  $\tilde{\Sigma}_2$  which is obtained by revolving the graph of the function  $\tilde{\rho}$  about the  $\zeta$ -axis.

When  $\tilde{\Omega}$  is "deformed" into  $\tilde{\Omega}_t$  by the speed V, the boundary  $\tilde{\Sigma}_2$  must remain a  $\zeta$ -axisymmetric graph. So we shall only consider speeds V of the following special form:

$$V(\rho,\zeta) = \begin{pmatrix} \omega(\rho,\zeta) \\ 0 \end{pmatrix}, \qquad \rho = \sqrt{\xi_1^2 + \xi_2^2}.$$
(5.4)

In the deformation the  $\zeta$ -axis and the boundary  $\tilde{\Sigma}_1$  both remain fixed, that is,

 $\omega(0,\zeta) = 0, \quad \forall \zeta \text{ in } [0,1]$  (5.5)

$$\omega(\rho, 0) = 0, \quad \exists \rho \text{ in } [0, 1]. \tag{5.6}$$

### 5.3. Derivation with Respect to the Shape

The scaled temperature  $y = y(\tilde{\Omega})$  is defined on the domain  $\tilde{\Omega}$ . In the "deformed" or "perturbed" domain  $\tilde{\Omega}_t$ , the new scaled temperature  $y_t = y(\tilde{\Omega}_t)$  is the solution of the following variational problem: to find  $y_t$  in  $H_0(\tilde{\Omega}_t)$  such that

$$\int_{\tilde{\Omega}_{t}} \left[ \tilde{L}^{2} \left( \frac{\partial y_{t}}{\partial \xi_{1}} \frac{\partial \psi}{\partial \xi_{1}} + \frac{\partial y_{t}}{\partial \xi_{2}} \frac{\partial \psi}{\partial \xi_{2}} \right) + \frac{\partial y_{t}}{\partial \zeta} \frac{\partial \psi}{\partial \zeta} \right] d\xi d\zeta = \int_{\tilde{\Sigma}_{1}} \psi d\tilde{\Sigma}_{1}$$
(5.7)

for all  $\psi$  in  $H_0(\tilde{\Omega}_t)$ , where

$$H_0(\tilde{\Omega}_t) = \{ \psi \in H^1(\tilde{\Omega}_t) | \psi = 0 \text{ on } \tilde{\Sigma}_3 \}.$$
(5.8)

Assume that the angles  $\theta_1$  and  $\theta_3$  (cf. sec. 3.4) are such that the solution  $y_t$  of (3.7) is itself smooth in a neighborhood of the curves  $C_1 = \tilde{\Sigma}_2 \cap \tilde{\Sigma}_1$  and  $C_3 = \tilde{\Sigma}_2 \cap \tilde{\Sigma}_3$ . Then the smoothness of  $y_t$  in the domain  $\tilde{\Omega}_t$  is given by its smoothness in an arbitrarily small neighborhood of the surface  $\tilde{\Sigma}_2$ . So further assume that the boundary  $\tilde{\Sigma}_2$  (that is, the functions  $\tilde{\rho}$  and  $\omega$ ) is sufficiently smooth in order to have the solution

$$y_t = y(\tilde{\Omega}_t) \in H^2(\tilde{\Omega}_t) \cap H_0(\tilde{\Omega}_t).$$
(5.9)

Any function in  $H^2(\tilde{\Omega}_t)$  may be extended to a function in  $H^2(\mathbb{R}^3) \to H^2(\mathbb{R}^3)$ . So look at  $y_t = y(\tilde{\Omega}_t)$  as the restriction to  $t \times \tilde{\Omega}_t$  of a smooth function Y defined on  $[0, t_1] \times \mathbb{R}^3$  (for some  $t_1 > 0$ ) such that for each t in  $[0, t_1]$  the map

$$x \to Y(t, x) : \mathbb{R}^3 \to \mathbb{R}$$

belongs to  $H^2(\mathbb{R}^3)$ . It is known from J.-P. Zolésio [12, 14], that Y may be chosen in such a way that

$$t \to Y(t, \cdot) : [0, t_1] \to H^1(\mathbb{R}^3)$$

is differentiable. Introduce the notation

$$Y' = \left. \frac{\partial Y}{\partial t} \right|_{t=0, x \in \tilde{\Omega}}.$$
(5.10)

Again from J. -P. Zolésio [12, 14], Y' is independent of the choice of the regular extension Y and we call Y' the *derivative* of  $y = y(\tilde{\Omega})$  for the domain  $\tilde{\Omega}$  in the *direction* V.

We shall now specify the boundary value problem of which Y' is a solution in the domain  $\tilde{\Omega}$ . This will require the introduction of elements of tangential differential calculus on the surface  $\tilde{\Sigma}_2$ .

#### 5.4. Tangential Calculus

Let  $\Sigma$  be the boundary of a smooth bounded domain  $\Omega$ , *n* be a normal vector field defined on  $\Sigma$ , *W* be a smooth vector field defined in a neighborhood of  $\Sigma$ , and *DW* be its Jacobian matrix.

It is known (cf. J.-P. Zolésio [12, 14]) that the expression

$$\operatorname{div} W - \langle DW, n, n \rangle \tag{5.11}$$

only depends on the restriction  $W|_{\Sigma}$  of the field W to the boundary  $\Sigma$ . Here " $\cdot$ " and " $\langle , \rangle$ " denote the product of a matrix by a vector and the inner product of two vectors, respectively.

So define the tangential divergence

$$\operatorname{div}_{T} W = \operatorname{div} W - \langle DW, n \rangle \text{ on } \Sigma.$$
(5.12)

It defines a differential operator on the surface  $\Sigma$ . This tangential divergence can also be defined by transposition in the following way. Given any  $\varphi$  in  $H^{s}(\Sigma)$ , s > 1, there exists  $\psi$  in  $H^{s+1/2}(\Omega)$  such that

$$\begin{cases} \psi = \varphi \quad \text{on } \Sigma \\ \frac{\partial \psi}{\partial n} = 0 \quad \text{on } \Sigma. \end{cases}$$
(5.13)

The gradient  $\nabla \psi$  on  $\Sigma$  is a tangent vector and is in fact independent of the choice of the function  $\psi$ . So define the *tangential gradient* of  $\varphi$  on the surface  $\Sigma$  as

$$\nabla_{\mathcal{T}} \varphi = \nabla \psi|_{\Sigma} \in H^{s-1}(\Sigma), \tag{5.14}$$

where  $\varphi$  is a given function only defined on  $\Sigma$ .

From J.-P. Zolésio [12], we have

$$\int_{\Sigma} \operatorname{div}_{T} W \varphi \, d\sigma = -\int_{\Sigma} \langle W, \nabla_{T} \varphi \rangle \, d\sigma - \int_{\Sigma} \langle H \varphi W, n \rangle \, d\sigma, \qquad (5.15)$$

where H is the "mean curvature" of the surface  $\Sigma$  (here we pick  $H = (k_1 + k_2)/2$ ,

where  $k_1$  and  $k_2$  are the upper and lower bounds of the normal curvatures on  $\Sigma$ , see, for instance, Do Carmo [1]).

### 5.5 Boundary-value Problem for Y'

Again from J.-P. Zolésio [12, 14], the shape derivative Y' of the scaled temperature  $y(\tilde{\Omega})$  in the direction V is the solution of the following variational problem: to find Y' in  $H_0(\tilde{\Omega})$  such that

$$\int_{\tilde{\Omega}} \left[ \tilde{L}^{2} \left( \frac{\partial Y'}{\partial \xi_{1}} \frac{\partial \varphi}{\partial \xi_{1}} + \frac{\partial Y'}{\partial \xi_{2}} \frac{\partial \varphi}{\partial \xi_{2}} \right) + \frac{\partial Y'}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] d\xi d\zeta$$
$$= -\int_{\tilde{\Sigma}_{2}} \left[ \tilde{L}^{2} \left( \frac{\partial y}{\partial \xi_{1}} \frac{\partial \varphi}{\partial \xi_{1}} + \frac{\partial y}{\partial \xi_{2}} \frac{\partial \varphi}{\partial \xi_{2}} \right) + \frac{\partial y}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] \langle V, \nu \rangle d\sigma \qquad (5.16)$$

for all  $\varphi$  in  $H^{3/2+\sigma}(\mathbb{R}^3)$ ,  $\sigma > 0$  arbitrary, such that  $\varphi = 0$  on  $\tilde{\Sigma}_3$ .

**Remark 5.1.** It is known from J. -P. Zolésio [14] that for  $s \ge 2$ 

$$y \in H^{s}(\tilde{\Omega}) \cap H_{0}(\tilde{\Omega}) \Rightarrow Y' \in H^{s-1}(\tilde{\Omega}) \cap H_{0}(\tilde{\Omega}).$$
 (5.17)

It follows from (5.16) that Y' is the solution of the following boundary-value problem:

$$\begin{array}{l} A Y' = 0 & \text{in } \tilde{\Omega} \\ Y' = 0 & \text{on } \tilde{\Sigma}_{3} \\ \\ \frac{\partial Y'}{\partial \nu_{A}} = 0 & \text{on } \tilde{\Sigma}_{1} \end{array} \right)$$

$$\begin{array}{l} \int_{\tilde{\Sigma}_{2}} \frac{\partial Y'}{\partial \nu_{A}} \varphi \, d\sigma = -\int_{\tilde{\Sigma}_{2}} \left[ \tilde{L}^{2} \left( \frac{\partial y}{\partial \xi_{1}} \frac{\partial \varphi}{\partial \xi_{1}} + \frac{\partial y}{\partial \xi_{2}} \frac{\partial \varphi}{\partial \xi_{2}} \right) + \frac{\partial y}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] \langle V, \nu \rangle \, d\sigma \\ \forall \varphi \in H^{3/2 + \sigma}(\tilde{\Omega}), \text{ such that } \exists \text{ a neighbourhood } V_{2} \text{ of } \tilde{\Sigma}_{2} \text{ where} \\ \varphi = 0 \text{ on } V_{2} \cap \tilde{\Sigma}_{3}. \end{array} \right)$$

$$(5.18)$$

For  $0 < \sigma < \frac{1}{2}$ , the closure in  $H^{3/2+\sigma}(\tilde{\Omega})$  of the set of functions

$$\left\{\varphi \in H^{3/2+\sigma}(\tilde{\Omega}) | \exists a \text{ neighborhood } V_2 \text{ of } \tilde{\Sigma}_2 \text{ such that } \varphi = 0 \text{ on } V_2 \cap \tilde{\Sigma}_3\right\}$$
(5.20)

coincides with

$$\Phi = \{\varphi \in H^{3/2+\sigma}(\tilde{\Omega}) | \varphi = 0 \text{ on } \tilde{\Sigma}_2 \cap \tilde{\Sigma}_3\}.$$
(5.21)

Recall that the injection

$$H^{3/2+\sigma}(\tilde{\Omega}) \to C^0(\overline{\tilde{\Omega}})$$
 (5.22)

is continuous and that the evaluation map

$$\varphi \to \varphi|_{\tilde{\Sigma}_2 \cap \tilde{\Sigma}_3} : C^0(\tilde{\Omega}) \to C^0(\tilde{\Sigma}_2 \cap \tilde{\Sigma}_3)$$
(5.23)

is also continuous. For  $\sigma$  in  $[0, \frac{1}{2}]$  there is no condition on the partial derivatives of  $\varphi$  on  $\tilde{\Sigma}_2 \cap \tilde{\Sigma}_3$ . As a result (5.19) is true for all  $\varphi$  in  $\Phi$  as defined in (5.21):

$$\int_{\tilde{\Sigma}_{2}} \frac{\partial Y'}{\partial \nu_{A}} \varphi \, d\sigma = -\int_{\tilde{\Sigma}_{2}} \left[ \tilde{L}^{2} \left( \frac{\partial y}{\partial \xi_{1}} \frac{\partial \varphi}{\partial \xi_{1}} + \frac{\partial y}{\partial \xi_{2}} \frac{\partial \varphi}{\partial \xi_{2}} \right) + \frac{\partial y}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] \langle V, \nu \rangle \, d\sigma$$

$$\forall \varphi \in H^{3/2+\sigma}(\tilde{\Omega}) \text{ such that } \varphi = 0 \text{ on } \tilde{\Sigma}_{2} \cap \tilde{\Sigma}_{3}, \qquad 0 < \sigma \leq \frac{1}{2}.$$

$$(5.24)$$

We now show that condition (5.24) is a Neuman condition on the boundary  $\tilde{\Sigma}_2$ . To see that we introduce the matrix

$$\mathscr{R} = \begin{bmatrix} \tilde{L}^2 & 0 & 0\\ 0 & \tilde{L}^2 & 0\\ 0 & 0 & 1 \end{bmatrix}$$
(5.25)

and perform the following computation on the expression

$$E = -\int_{\tilde{\Sigma}_2} \left[ \tilde{L}^2 \left( \frac{\partial y}{\partial \xi_1} \frac{\partial \varphi}{\partial \xi_1} + \frac{\partial y}{\partial \xi_2} \frac{\partial \varphi}{\partial \xi_2} \right) + \frac{\partial y}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] \langle V, \nu \rangle \, d\sigma.$$
(5.26)

By rewriting

$$E = -\int_{\tilde{\Sigma}_{2}} \langle \langle V, \nu \rangle \mathscr{R} \cdot \nabla y, \nabla \varphi \rangle \, d\sigma, \qquad (5.27)$$

where  $\mathscr{Q} \cdot \nabla y$  is the vector

$$\left(\tilde{L}^2 \frac{\partial y}{\partial \xi_1}, \tilde{L}^2 \frac{\partial y}{\partial \xi_2}, \frac{\partial y}{\partial \zeta}\right)$$
(5.28)

Define the field

$$W = \langle V, \mathbf{v} \rangle \mathcal{C} \cdot \nabla y \tag{5.29}$$

and rewrite the right-hand side of (5.27) as

$$E = -\int_{\tilde{\Sigma}_{2}} \langle W, \nabla \varphi \rangle \, d\sigma$$
  
=  $-\int_{\tilde{\Sigma}_{2}} \langle W, \nabla_{T} \varphi \rangle \, d\sigma - \int_{\tilde{\Sigma}_{2}} \langle W, \frac{\partial \varphi}{\partial \nu} \nu \rangle \, d\sigma.$  (5.30)

The second term on the right-hand side of (5.30) can be rewritten in the form

$$\int_{\tilde{\Sigma}_2} \frac{\partial \varphi}{\partial \nu} \langle W, \nu \rangle \, d\sigma. \tag{5.31}$$

But  $\langle W, \nu \rangle$  is zero and the integral (5.31) is also equal to zero: by definition of W

$$\langle W, \nu \rangle = \langle V, \nu \rangle \langle \nu, \mathcal{Q} \cdot \nabla y \rangle$$
  
=  $\langle V, \nu \rangle \langle \mathcal{Q} \cdot \nu, \nabla y \rangle = \langle V, \nu \rangle \frac{\partial y}{\partial \nu_A} = 0$ 

on  $\tilde{\Sigma}_2$ . As for the first term on the right-hand side of (5.30), it can be evaluated with the help of identity (5.15):

$$\int_{\tilde{\Sigma}} \operatorname{div}_{T} W \varphi \, d\sigma = -\int_{\tilde{\Sigma}} \langle W, \nabla_{T} \varphi \rangle \, d\sigma - \int_{\tilde{\Sigma}} H \varphi \langle W, \nu \rangle \, d\sigma \tag{5.32}$$

with the vector field W chosen as in (5.29). On  $\tilde{\Sigma}_1 \cup \tilde{\Sigma}_3$ , W is zero since  $\langle V, \nu \rangle$  is zero by our choice of the speed V. So we can substitute  $\tilde{\Sigma}_2$  for  $\tilde{\Sigma}$  in identity (5.32). Moreover, the second term in that identity is zero over  $\tilde{\Sigma}_2$  since we have already shown that  $\langle W, \nu \rangle$  is zero on  $\tilde{\Sigma}_2$ . Finally we obtain the following chain of identities:

$$\int_{\tilde{\Sigma}_{2}} \operatorname{div}_{T} W \varphi|_{\tilde{\Sigma}_{2}} d\sigma = -\int_{\tilde{\Sigma}_{2}} \langle W, \nabla_{T} \varphi \rangle d\sigma = -\int_{\tilde{\Sigma}_{2}} \langle W, \nabla \varphi \rangle d\sigma$$
(5.33)

for all  $\varphi$  in  $H^{3/2+\sigma}(\tilde{\Omega})$ . Notice that the above identities do not require that  $\varphi$  be zero on  $\tilde{\Sigma}_3$ . From identities (5.19) and (5.27)

$$\int_{\tilde{\Sigma}_{2}} \frac{\partial Y'}{\partial \nu_{A}} \varphi \, d\sigma = -\int_{\tilde{\Sigma}_{2}} \langle W, \nabla \varphi \rangle \, d\sigma, \forall \varphi \in H_{0}(\tilde{\Omega}) \cap H^{3/2+\sigma}(\tilde{\Omega}).$$
(5.34)

Combining identities (5.33) and (5.34)

$$\int_{\tilde{\Sigma}_{2}} \frac{\partial Y'}{\partial \nu_{\mathcal{A}}} \varphi \bigg|_{\tilde{\Sigma}_{2}} d\sigma = \int_{\tilde{\Sigma}_{2}} \operatorname{div}_{T} W \varphi \bigg|_{\tilde{\Sigma}_{2}} d\sigma, \forall \varphi \in H_{0}(\tilde{\Omega}) \cap H^{3/2+\sigma}(\tilde{\Omega}).$$
(5.35)

For each h in  $\mathfrak{P}(\tilde{\Sigma}_2)$ , there exists  $\varphi$  in  $H_0(\tilde{\Omega}) \cap H^{3/2+\sigma}(\tilde{\Omega})$  such that

$$A\varphi = 0, \qquad \varphi \Big|_{\tilde{\Sigma}_3} = 0, \frac{\partial \varphi}{\partial \nu_A} \Big|_{\tilde{\Sigma}_2} = h, \qquad \frac{\partial \varphi}{\partial \nu_A} \Big|_{\tilde{\Sigma}_1} = 0.$$

As a result

$$\int_{\tilde{\Sigma}_{2}} \frac{\partial Y'}{\partial \nu_{A}} h \, d\sigma = \int_{\tilde{\Sigma}_{2}} \operatorname{div}_{T} W h \, d\sigma, \quad \forall h \in \mathfrak{N}(\tilde{\Sigma}_{2}).$$
(5.36)

But for y in  $H^{2+\epsilon}(\tilde{\Omega})$ ,  $\epsilon \ge 0$ ,  $W \in H^{1+\epsilon}(\tilde{\Omega})$  and  $\operatorname{div}_T W \in H^{-1/2+\epsilon}(\tilde{\Sigma}_2)$  (use the fact that y is harmonic to define the trace  $\operatorname{div}_T W$ ). So the right-hand side of (5.36) is continuous for the  $H^{1/2-\epsilon}(\tilde{\Sigma}_2)$ -topology. By density of  $\mathfrak{N}(\tilde{\Sigma}_2)$  in  $H^{1/2-\epsilon}(\tilde{\Sigma}_2)$  for  $\epsilon \ge 0$ , identity (5.36) holds for all h in  $H^{1/2-\epsilon}(\tilde{\Sigma}_2)$  and necessarily

$$\frac{\partial Y'}{\partial \nu_A} = \operatorname{div}_T W = \operatorname{div}_T (\mathscr{Q} \cdot \nabla y \langle V, \nu \rangle) \text{ in } H^{-1/2 + \varepsilon} (\tilde{\Sigma}_2).$$
(5.37)

We summarize the following useful relations:

$$\begin{array}{l}
\nu_{A} = \mathscr{Q} \cdot \nu \quad (\text{conormal associated with the op-} \\
\text{erator } A \text{ of (3.11) and the normal} \\
\text{vector } \nu) \\
\frac{\partial \varphi}{\partial \nu_{A}} = \langle \mathscr{Q} \cdot \nu, \nabla \varphi \rangle = \langle \nu_{A}, \nabla \varphi \rangle \\
A\varphi = -\operatorname{div}(\mathscr{Q} \cdot \nabla \varphi).
\end{array}$$
(5.38)

If the cotangential divergence (with respect to the operator A) is defined as

$$\operatorname{div}_{T_{4}}(\vec{E}) = \operatorname{div}_{T}(\mathscr{Q}.\vec{E}), \tag{5.39}$$

where  $\vec{E}$  is a vector function from  $\tilde{\Sigma}$  into  $\mathbb{R}^3$ , then (5.37) is equivalent to

$$\frac{\partial Y'}{\partial \nu_A} = \operatorname{div}_{T_A}(\nabla y \langle V, \nu \rangle).$$
(5.40)

### 6. Eulerian Derivative of the Cost and the Constraint

Recall from sec. 3.1 that the cost function to be minimized is of the form

$$J(\tilde{L},\tilde{\rho}) = \tilde{L}J(\tilde{\Omega}), \tag{6.1}$$

where

$$J(\tilde{\Omega}) = \pi \int_0^1 \tilde{\rho}(\zeta)^2 d\zeta \bigg( = \int_{\tilde{\Omega}} d\xi_1 d\xi_2 d\zeta, \text{ the volume of } \tilde{\Omega} \bigg).$$
(6.2)

Following J.-P. Zolésio [12, 14], define (whenever it exists) the *Eulerian derivative*  $dJ(\tilde{\Omega}; V)$  (or shape derivative) at  $\tilde{\Omega}$  in the direction of the field V as

$$dJ(\tilde{\Omega}; V) = \lim_{t \searrow 0} \left[ J(\tilde{\Omega}_t) - J(\tilde{\Omega}) \right] / t.$$
(6.3)

# 6.1. Shape Derivative of the Volume

From J.-P. Zolésio [12] the shape derivative of  $J(\tilde{\Omega})$  is

$$dJ(\tilde{\Omega};V) = \int_{\tilde{\Sigma}} \langle V, \nu \rangle \, d\sigma, \qquad (6.4)$$

where  $\langle V, \nu \rangle$  is the normal component of the field V on  $\tilde{\Sigma}$ . It is equal to zero on  $\tilde{\Sigma}_1$  and  $\tilde{\Sigma}_3$  and

$$\langle V, \nu \rangle = \omega \nu_r \text{ on } \tilde{\Sigma}_2.$$
 (6.5)

Finally

$$dJ(\tilde{L},\tilde{\rho};V) = \tilde{L} dJ(\tilde{\Omega};V) = \tilde{L} \int_{\tilde{\Sigma}} \omega \nu_r d\sigma$$
  
=  $\tilde{L} \int_0^1 2\pi \omega(\tilde{\rho}(\zeta),\zeta) \tilde{\rho}(\zeta) d\zeta.$  (6.6)

6.2. Shape Derivative of the Constraint Functional f

Recall that

$$f(\tilde{L},\tilde{\rho}) = \int_{\tilde{\Sigma}_3} \left[ \frac{\partial y}{\partial \zeta} + q \right]^- d\sigma = 2\pi \int_0^{\tilde{\rho}(1)} \left[ \frac{\partial y}{\partial \zeta}(\rho, 1) + q \right]^- \rho \, d\rho \tag{6.7}$$

and that  $Y(t, \rho, \zeta)$  is a smooth extension of  $y = y(\tilde{\Omega})$ . As a result for t > 0,

$$f(\tilde{L},\tilde{\rho}_t) = 2\pi \int_0^{\tilde{\rho}_t(1)} \left[ \frac{\partial y}{\partial \zeta}(t;\rho,1) + q \right]^- \rho \, d\rho.$$
(6.8)

We consider two cases:

1) the constraint is saturated everywhere on the boundary  $\tilde{\Sigma}_3$ ,

$$-\left[\frac{\partial y}{\partial \zeta}(\rho,1)+q\right] > 0, \quad \forall \rho, \quad 0 \leq \rho < \tilde{\rho}(1),$$
(6.9)

2) the constraint is not saturated on a subset of the boundary  $\tilde{\Sigma}_3$  of non-zero measure.

In both cases we know that

$$\int_{\tilde{\Sigma}} \frac{\partial y}{\partial \nu_A} d\sigma = 0 \Rightarrow -\int_{\tilde{\Sigma}_3} \frac{\partial y}{\partial \nu_A} d\nu = \int_{\tilde{\Sigma}_1} \frac{\partial y}{\partial \nu_A} d\sigma = \pi.$$
(6.10)

As a result

$$-\frac{\partial y}{\partial \nu_{A}}(\rho,1) \leqslant q, \qquad \forall 0 \leqslant \rho \leqslant \tilde{\rho}(1) \Rightarrow \tilde{\rho}(1) \geqslant \sqrt{1/q}$$
(6.11)

$$-\frac{\partial y}{\partial \nu_{A}}(\rho,1) > q, \qquad \forall 0 \le \rho < \tilde{\rho}(1) \Rightarrow \tilde{\rho}(1) < \sqrt{1/q}.$$
(6.12)

The condition

$$\tilde{\rho}(1) \ge \sqrt{1/q} \tag{6.13}$$

is *necessary* in order to satisfy the constraint on the boundary  $\tilde{\Sigma}_3$ . In case 1 condition (6.13) will be violated. In order to eliminate that situation we make the following hypothesis.

*Hypothesis 6.1.* The shape function  $\tilde{\rho}$  will satisfy the following conditions

$$\tilde{\rho}(0) = 1 \quad \text{and} \quad \tilde{\rho}(1) \ge \sqrt{1/q} \,.$$
(6.14)

Remark 6.1. In case 1

$$f(\tilde{L},\tilde{\rho}) = \int_{\tilde{\Sigma}_3} \left[ \frac{\partial y}{\partial \nu_A} + q \right]^- d\sigma = -\int_{\tilde{\Sigma}_3} \frac{\partial y}{\partial \nu_A} d\sigma - q \pi \tilde{\rho}(1)^2$$
$$= \pi \left( 1 - q \tilde{\rho}(1)^2 \right)$$
(6.15)

and

$$df(\tilde{L},\tilde{\rho};V) = -q2\pi\tilde{\rho}(1)\omega(\tilde{\rho}(1),1), \qquad (6.16)$$

where  $V = (\omega, 0)$  is the speed of the points in the domain. In this situation the gradient says that the point  $(1, \tilde{\rho}(1))$  has to be moved away from the  $\zeta$ -axis, but *it does not say* what should be done with the remainder,  $\langle (\zeta, \tilde{\rho}(\zeta) : 0 \leq \zeta < 1 \rangle$ , of the boundary  $\tilde{\Sigma}_2$ . Hypothesis 6.1 will ensure that case 1 and this situation do not occur.

In the second case, we introduce the following hypothesis.

*Hypothesis 6.2.* There exists a neighborhood N of the curve  $C_3 = \tilde{\Sigma}_2 \cap \tilde{\Sigma}_3$  such that

$$\left[\frac{\partial y}{\partial \zeta}(\sigma) + q\right] = 0, \quad \forall \sigma \in \tilde{\Sigma}_3 \cap N.$$
(6.17)

Under that hypothesis the integrand of the constraint (6.7) is zero in the neighborhood of  $\tilde{\rho}(1)$ . So there will be no boundary term on  $C_3$  and the

directional derivative of  $f(\tilde{L}, \tilde{\rho}_t; V)$  at t = 0 will be given by

$$df(\tilde{L},\tilde{\rho};V) = \frac{d}{dt} \left\{ 2\pi \int_0^{\tilde{\rho}(1)} \left[ \frac{\partial y}{\partial \zeta}(t;\rho,1) + q \right]^- \rho \, d\rho \right\}_{t=0}, \tag{6.18}$$

where d/dt denotes the right-hand side derivative. By using general results in J.-P. Zolésio [12, 14], we know that expression (6.13) is equal to

$$-2\pi \int_{0}^{\tilde{\rho}(1)} \frac{\partial Y'}{\partial \zeta}(\rho, 1) \chi_{+}(\rho) \rho \, d\rho + 2\pi \int_{0}^{\tilde{\rho}(1)} \left[\frac{\partial Y'}{\partial \zeta}(\rho, 1)\right]^{-} \chi_{0}(\rho) \rho \, d\rho, \quad (6.19)$$

where

$$\chi_{+}(\rho)(\operatorname{resp.} \chi_{0}(\rho)) = \begin{cases} 1, & \text{if } \frac{\partial y}{\partial \zeta}(\rho, 1) + q < 0 \text{ (resp. = 0)} \\ 0, & \text{otherwise.} \end{cases}$$
(6.20)

We introduce another hypothesis in order to obtain the linearity of the derivative  $df(\tilde{\Omega}; V)$  with respect to V.

*Hypothesis 6.3.* For physical reasons, we shall assume that when 0 < q < 1 the subset of all  $\rho, 0 \le \rho \le \tilde{\rho}(1)$ , such that

$$-\frac{\partial y}{\partial \nu_{A}}(\rho,1) = q \tag{6.21}$$

has zero measure.

So the characteristic function  $\chi_0$  is zero almost everywhere in  $[0, \tilde{\rho}(1)]$  and the second term in (6.19) is zero. Under Hypotheses 6.1, 6.2, and 6.3

$$df(\tilde{L},\tilde{\rho};V) = -2\pi \int_0^{\tilde{\rho}(1)} \frac{\partial Y'}{\partial \zeta}(\rho,1) \chi_+(\rho) \rho \, d\rho.$$
(6.22)

Hypotheses 6.2 and 6.3 can be combined into a single stronger hypothesis.

Hypothesis 6.4. The function

$$\rho \to -\frac{\partial y}{\partial \zeta}(\rho, 1) : [0, \tilde{\rho}(1)] \to \mathbb{R}$$
(6.23)

is monotone strictly increasing.

#### 6.3. Derivative of the Penalized Cost

Under hypotheses 6.1, 6.2, and 6.3 and in view of (4.3), (6.6) and (6.22), the Eulerian derivative of the penalized cost is given by the following expression:

$$\frac{1}{2\pi}dJ_{\varepsilon}(\tilde{L},\tilde{\rho};V) = \tilde{L}\int_{0}^{1}\omega(\tilde{\rho}(\zeta),\zeta)\tilde{\rho}(\zeta)\,d\zeta - \frac{1}{\varepsilon}\int_{0}^{\tilde{\rho}(1)}\frac{\partial Y'}{\partial\zeta}(\rho,1)\chi_{+}(\rho)\rho\,d\rho.$$
(6.24)

# 7. Derivative with Respect to the Parameter $\tilde{L}$

7.1. Derivative of the Cost Function

Trivially

$$\frac{\partial J}{\partial \tilde{L}}(\tilde{L},\tilde{\rho}) = J(\tilde{\Omega}) = \pi \int_0^1 \tilde{\rho}(\zeta)^2 d\zeta$$
(7.1)

# 7.2. Derivative of the Scaled Temperature

Let y be the solution of problem (5.7) for t = 0. By the implicit function theorem (cf. J. -P. Zolésio [14], G. Chavent [3]) we know that the maps  $\tilde{L} \to y : \mathbb{R}_+ \to H^2(\tilde{\Omega})$  is differentiable. So  $y_L = \frac{\partial y}{\partial \tilde{L}}$  is the solution of the following variational problem: to find  $y_L$  in  $H_0(\tilde{\Omega})$  such that

$$\int_{\tilde{\Omega}} \left[ \tilde{L}^{2} \left( \frac{\partial y_{L}}{\partial \xi_{1}} \frac{\partial \psi}{\partial \xi_{1}} + \frac{\partial y_{L}}{\partial \xi_{2}} \frac{\partial \psi}{\partial \xi_{2}} \right) + \frac{\partial y_{L}}{\partial \zeta} \frac{\partial \psi}{\partial \zeta} \right] d\xi d\zeta$$
  
=  $-2 \tilde{L} \int_{\tilde{\Omega}} \left( \frac{\partial y}{\partial \xi_{1}} \frac{\partial \psi}{\partial \xi_{1}} + \frac{\partial y}{\partial \xi_{2}} \frac{\partial \psi}{\partial \xi_{2}} \right) d\xi d\zeta$  (7.2)

for all  $\psi$  in  $H_0(\tilde{\Omega})$ .

We can further specify the boundary-value problem of which  $y_L$  is a solution:

$$Ay_{L} = -\left[\tilde{L}^{2}\left(\frac{\partial^{2}y_{L}}{\partial\xi_{1}^{2}} + \frac{\partial^{2}y_{L}}{\partial\xi_{2}^{2}}\right) + \frac{\partial^{2}y_{L}}{\partial\xi^{2}}\right] = 2\tilde{L}\left(\frac{\partial^{2}y}{\partial\xi_{1}^{2}} + \frac{\partial^{2}y}{\partial\xi_{2}^{2}}\right) \text{ in }\tilde{\Omega},$$
(7.3)  
$$\frac{\partial y_{L}}{\partial\nu_{A}} = \tilde{L}^{2}\left(\nu_{1}\frac{\partial y_{L}}{\partial\xi_{1}} + \nu_{2}\frac{\partial y_{L}}{\partial\xi_{2}}\right) + \frac{\partial y_{L}}{\partial\xi}$$

$$= -2\tilde{L}\left(\nu_1\frac{\partial y}{\partial\xi_1} + \nu_2\frac{\partial y}{\partial\xi_2}\right) \text{ on } \tilde{\Sigma}_1 \cup \tilde{\Sigma}_2$$
(7.4)

Using equations (3.11)–(3.12) for y the two equations (7.3)–(7.4) are equivalent to

$$Ay_L = -\frac{2}{\tilde{L}} \frac{\partial^2 y}{\partial \zeta^2}$$
(7.3a)

$$\frac{\partial y_L}{\partial \nu_A}\Big|_{\tilde{\Sigma}_1} = 1, \frac{\partial y_L}{\partial \nu_A}\Big|_{\tilde{\Sigma}_2} = \frac{2}{\tilde{L}} \nu_{\zeta} \frac{\partial y}{\partial \zeta}, y_L\Big|_{\tilde{\Sigma}_3} = 0.$$
(7.4a)

For the condition on  $\tilde{\Sigma}_1$  we have used the fact that on  $\tilde{\Sigma}_1$ 

$$\nu = (0, 0, -1) \Rightarrow \nu_1 = \nu_2 = 0.$$
 (7.5)

**Remark 7.1.** In view of the boundary conditions (7.4a),  $y_L$  has the same smoothness as y provided the boundary  $\tilde{\Sigma}_2$  is sufficiently smooth.

#### 7.3. Derivative of the Constraint Functional

As in section 6.2 we compute the right-hand side derivative:

$$\frac{df}{dL}(\tilde{L},\tilde{\rho}) = -2\pi \int_{0}^{\tilde{\rho}(1)} \frac{\partial y_{L}}{\partial \zeta}(\rho,1) \chi_{+}(\rho) \rho \, d\rho + 2\pi \int_{0}^{\tilde{\rho}(1)} \left[\frac{\partial y_{L}}{\partial \zeta}(\rho,1) + q\right]^{-} \chi_{0}(\rho) \rho \, d\rho,$$
(7.6)

where  $\chi_0$  and  $\chi_+$  are as defined in (6.11). For physical reasons we have already assumed that  $\chi_0$  is zero almost everywhere (cf. Hypothesis 6.3). As a result

$$\frac{d}{dL}f(\tilde{L},\tilde{\rho}) = -2\pi \int_0^{\tilde{\rho}(1)} \frac{\partial y_L}{\partial \zeta}(\rho,1) \chi_+(\rho) \rho \, d\rho.$$
(7.7)

7.4. Derivative of the Penalized Cost

In view of (4.3), (7.1) and (7.6)

$$\frac{1}{2\pi} \frac{\partial J_{\varepsilon}}{\partial \tilde{L}} (\tilde{L}, \tilde{\rho}) = \frac{1}{2} \int_0^1 \tilde{\rho}^2(\zeta) \, d\zeta - \frac{1}{\varepsilon} \int_0^{\tilde{\rho}(1)} \frac{\partial Y_L}{\partial \zeta} (\rho, 1) \chi_+(\rho) \rho \, d\rho.$$
(7.8)

### 8. Gradient Computations

Recall expressions (6.13) for the derivative of the penalized cost with respect to the shape of the volume  $\tilde{\Omega}$  and expression (7.8) for the derivative of the penalized cost with respect to the parameter  $\tilde{L}$ . They both involve a term of the form

$$\lambda(\psi) = -\int_{\tilde{\Sigma}_3} \frac{\partial \psi}{\partial \nu_A} \chi_+ \, d\sigma = -2\pi \int_0^{\tilde{\rho}(1)} \frac{\partial \psi}{\partial \zeta}(\rho, 1) \chi_+(\rho) \rho \, d\rho \tag{8.1}$$

with  $\psi = Y'$  and  $\psi = y_L$ , respectively.

Following a standard procedure we shall now introduce the solution p of an appropriate adjoint system and express  $\lambda(Y')$  and  $\lambda(y_L)$  as a function of y, p and the speed V. The following bilinear forms will be useful in the forthcoming discussion:

$$a(\psi,\varphi) = \int_{\hat{\Omega}} \left[ \tilde{L}^2 \left( \frac{\partial \psi}{\partial \xi_1} \frac{\partial \varphi}{\partial \xi_1} + \frac{\partial \psi}{\partial \xi_2} \frac{\partial \varphi}{\partial \xi_2} \right) + \frac{\partial \psi}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] d\xi \, d\zeta \tag{8.2}$$
$$(d\xi = d\xi_1 d\xi_2),$$

$$b(\psi,\varphi) = -2\tilde{L} \int_{\tilde{\Omega}} \left( \frac{\partial \psi}{\partial \xi_1} \frac{\partial \varphi}{\partial \xi_1} + \frac{\partial \psi}{\partial \xi_2} \frac{\partial \varphi}{\partial \xi_2} \right) d\xi d\zeta$$
(8.3)

and

$$B(\psi,\varphi) = -\int_{\tilde{\Sigma}_2} \left[ \tilde{L}^2 \left( \frac{\partial \psi}{\partial \xi_1} \frac{\partial \varphi}{\partial \xi_1} + \frac{\partial \psi}{\partial \xi_2} \frac{\partial \varphi}{\partial \xi_2} \right) + \frac{\partial \psi}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] \langle V, \nu \rangle \, d\sigma \tag{8.4}$$

### 8.1. Adjoint Equation

Let *p* be the solution of the system

$$Ap = -\left[\tilde{L}^{2}\left(\frac{\partial^{2}p}{\partial\xi_{1}^{2}} + \frac{\partial^{2}p}{\partial\xi_{2}^{2}}\right) + \frac{\partial^{2}p}{\partial\xi^{2}}\right] = 0 \text{ in } \tilde{\Omega}$$

$$p = \chi_{+} \text{ on } \tilde{\Sigma}_{3}$$

$$\frac{\partial p}{\partial\nu_{A}} = 0 \text{ on } \tilde{\Sigma}_{1} \cup \tilde{\Sigma}_{2}.$$
(8.5)

The function  $\chi_+$  is the characteristic function of a disk centered in r = 0 on the surface  $\tilde{\Sigma}_3$ . So we only have

$$\chi_{+} \in H^{1/2-\varepsilon}(\tilde{\Sigma}_{3}), \quad \forall \varepsilon > 0.$$
(8.6)

As a result p belongs to  $H^{1-\epsilon}(\tilde{\Omega})$  and even to  $D_{A}^{1-\epsilon}(\tilde{\Omega})$ , where for 0 < s < 2,

$$D_{A}^{s}(\tilde{\Omega}) = \left\{ v \in H^{s}(\tilde{\Omega}) : Av \in \Xi E^{s-2}(\tilde{\Omega}) \right\}$$
  

$$\Xi^{s-2}(\tilde{\Omega}) = \left( \Xi^{2-s}(\tilde{\Omega}) \right)'(' = \text{topological dual})$$
  

$$\Xi^{2-s}(\tilde{\Omega}) = \left\{ v \in \mathfrak{D}'(\tilde{\Omega})' : d^{2-s}v \in H_{0}^{2-s}(\tilde{\Omega}) \right\}$$
(8.7)

and  $d(x, \tilde{\Sigma})$  is the distance from x to the boundary  $\tilde{\Sigma}$  (cf. Lions and Magenes [11, Vol. 1, pp. 183, 199]).

We shall express  $\lambda(\psi)$  as a function of p. For all  $\varphi$  and  $\psi$  in  $H^2(\tilde{\Omega})$  the following identity holds

$$(A\psi,\varphi)_{L^{2}(\tilde{\Omega})} + \left(\frac{\partial\psi}{\partial\nu_{A}},\varphi\right)_{L^{2}(\tilde{\Sigma})} = (\psi,A\varphi)_{L^{2}(\tilde{\Omega})} + \left(\psi,\frac{\partial\varphi}{\partial\nu_{A}}\right)_{L^{2}(\tilde{\Sigma})}.$$
(8.8)

When  $\varphi$  and  $\psi$  are such that

$$A\varphi = 0, \frac{\partial \varphi}{\partial \nu_{A}} = 0 \text{ on } \tilde{\Sigma}_{1} \cup \tilde{\Sigma}_{2} \text{ and } \psi|_{\tilde{\Sigma}_{3}} = 0, \qquad (8.9)$$

the right hand side of identity (8.8) is zero and the following new identity holds

$$-\left(\frac{\partial\psi}{\partial\nu_{A}},\varphi\right)_{L^{2}(\tilde{\Sigma}_{3})} = (A\psi,\varphi)_{L^{2}(\tilde{\Omega})} + \left(\frac{\partial\psi}{\partial\nu_{A}},\varphi\right)_{L^{2}(\tilde{\Sigma}_{1}\cup\tilde{\Sigma}_{2})}.$$
(8.10)

By continuity the right-hand side of (8.10) makes sense for all  $\varphi$  in  $H^{1-\epsilon}(\tilde{\Omega})$ ,  $\epsilon < 1/2$ , for which identities (8.9) are verified. In particular one can pick  $\varphi = p$ , the solution of the adjoint system (8.5). Finally

$$\lambda(\psi) = -\left(\frac{\partial\psi}{\partial\nu_{A}}, \chi_{+}\right)_{L^{2}(\tilde{\Sigma}_{3})} = (A\psi, p)_{L^{2}(\tilde{\Omega})} + \left(\frac{\partial\psi}{\partial\nu_{A}}, p\right)_{L^{2}(\tilde{\Sigma}_{1}\cup\tilde{\Sigma}_{2})}$$
(8.11)

for all  $\psi$  in  $H^2(\tilde{\Omega}) \cap H_0(\tilde{\Omega})$ .

It is interesting that the right-hand side of expression (8.11) is continuous for  $\psi$  in  $D_A^{1+\epsilon}(\tilde{\Omega}) \cap H_0(\tilde{\Omega})$  since *p* belongs to  $H^{1-\epsilon}(\tilde{\Omega})$  and  $p|_{\tilde{\Sigma}_1 \cup \tilde{\Sigma}_2}$  to  $H^{1/2-\epsilon}(\tilde{\Sigma}_1 \cup \tilde{\Sigma}_2)$ . However in that case expression (8.11) must be written with duality products:

$$\lambda(\psi) = -\left\langle \frac{\partial \psi}{\partial \nu_{A}}, \chi_{+} \right\rangle_{H^{1/2-\epsilon}(\tilde{\Sigma}_{3})} = \langle A\psi, p \rangle_{\Xi^{1-\epsilon}(\tilde{\Omega})} + \left\langle \frac{\partial \psi}{\partial \nu_{A}}, p \right\rangle_{H^{1/2-\epsilon}(\tilde{\Sigma}_{1} \cup \tilde{\Sigma}_{2})}$$
(8.12)

for all  $\psi$  in  $D_A^{1+\epsilon}(\tilde{\Omega}) \cap H_0(\tilde{\Omega})$ .

It is useful to summarize the above discussion in the following proposition.

**Proposition 8.1.** For all  $\varepsilon$ ,  $0 < \varepsilon < \frac{1}{2}$ , and all  $\psi$  in  $D_A^{1+\varepsilon}(\tilde{\Omega}) \cap H_0(\tilde{\Omega})$ , expression (8.12) holds. If, in addition,  $\psi$  belongs to  $H^2(\tilde{\Omega})$ , expression (8.12) reduces to expression (8.11).

**Remark 8.1.** Equation (8.12) suggests that p could be obtained by transposition of the linear map

$$\psi \to A\psi, \left. \frac{\partial \psi}{\partial \nu_A} \right|_{\tilde{\Sigma}_1 \cup \tilde{\Sigma}_2} : D_A^{1+\epsilon}(\tilde{\Omega}) \cap H_0(\tilde{\Omega}) \to \Xi^{-1+\epsilon}(\tilde{\Omega}) \times H^{1/2+\epsilon}(\tilde{\Sigma}_1 \cup \tilde{\Sigma}_2)$$
(8.13)

provided that it is an isomorphism and that the various spaces  $D_A^{1+\epsilon}$ ,  $\Xi^{-1+\epsilon}$  and  $H^{1/2+\epsilon}$  are well-defined. This would require an adaptation of a result in Lions and Magenes [11, Vol. 1, Thm. 7.4]. With such spaces, the linear map

$$\psi \to \lambda(\psi) = -\left\langle \frac{\partial \psi}{\partial \nu_{\mathcal{A}}}, \chi_{+} \right\rangle_{H^{1/2-\epsilon}(\tilde{\Sigma}_{1} \cup \tilde{\Sigma}_{2})} : D_{\mathcal{A}}^{1+\epsilon}(\tilde{\Omega}) \cap H_{0}(\tilde{\Omega}) \to \mathbb{R}$$
(8.14)

would be continuous. So by the method of transposition one could claim that there exists a unique pair

$$p_0, p_1 \in \Xi^{1-\epsilon}(\tilde{\Omega}) \times H^{-1/2-\epsilon}(\tilde{\Sigma}_1 \cup \tilde{\Sigma}_2)$$
(8.15)

such that

$$\langle A\psi, p_0 \rangle_{\Xi^{1-\epsilon}(\tilde{\Omega})} + \left\langle \frac{\partial \psi}{\partial \nu_A}, p_1 \right\rangle_{H^{1/2-\epsilon}(\tilde{\Sigma}_1 \cup \tilde{\Sigma}_2)} = \lambda(\psi)$$
 (8.16)

for all  $\psi$  in  $D_A^{1+\epsilon}(\tilde{\Omega}) \cap H_0(\tilde{\Omega})$ . Now by comparison of (8.12) and (8.16) and uniqueness of the solution to (8.16), we should formally obtain the identities

$$p_0 = p \text{ in } \Omega \quad \text{and} \quad p_1 = p|_{\tilde{\Sigma}_1 \cup \tilde{\Sigma}_2}.$$
 (8.17)

**Remark 8.2.** When  $\chi_+ = 1$  on  $\tilde{\Sigma}_3$ , the solution p of (8.5) is identically equal to 1 in  $\tilde{\Omega}$ , the closure of  $\tilde{\Omega}$ .

### 8.2. Computation of $\lambda(Y')$ .

Recall from section 5.5 that Y' is the solution of the boundary value problem

$$\begin{cases} A Y' = 0 \text{ in } \tilde{\Omega} \\ \frac{\partial Y'}{\partial \nu_{A}} \Big|_{\tilde{\Sigma}_{1}} = 0, \ \frac{\partial Y'}{\partial \nu_{A}} \Big|_{\tilde{\Sigma}_{2}} = \operatorname{div}_{T}(\mathscr{Q} \cdot \nabla y \langle V, \nu \rangle), \qquad Y'|_{\tilde{\Sigma}_{3}=0}. \end{cases}$$
(8.18)

Also from Remark 5.1 the smoothness of Y' is related to the one of y as follows

$$y \in H^{s}(\tilde{\Omega}) \Rightarrow Y' \in H^{s-1}(\tilde{\Omega}), \quad s \ge 2,$$

provided that the boundary  $\tilde{\Sigma}_2$  is smooth enough and that the angles  $\theta_1$  and  $\theta_3$  are small enough.

So we assume for the moment that y belongs to  $H^3(\tilde{\Omega})$  which means that Y' belongs to  $H^2(\tilde{\Omega}) \cap H_0(\tilde{\Omega})$ . Then set  $\psi = Y'$  in expression (8.11):

$$\lambda(Y') = (AY', p)_{L^{2}(\tilde{\Omega})} + \left(\frac{\partial Y'}{\partial \nu_{A}}, p\right)_{L^{2}(\tilde{\Sigma}_{1} \cup \tilde{\Sigma}_{2})}.$$
(8.19)

In view of (8.18), identity (8.19) reduces to

$$\lambda(Y') = \left(\frac{\partial Y'}{\partial \nu_A}, p\right)_{L^2(\tilde{\Sigma}_2)}.$$
(8.20)

We have seen in section 5.5 that (cf. (5.31))

$$\frac{\partial Y'}{\partial \nu_A} = \operatorname{div}_T W, \qquad W = \mathcal{R}. \nabla y \langle V, \nu \rangle$$
(8.21)

and that for all  $\varphi$  in  $H^{3/2+\sigma}(\tilde{\Omega})$  (cf. (5.27))

$$\int_{\tilde{\Sigma}_{2}} \operatorname{div}_{T} W \varphi|_{\tilde{\Sigma}_{2}} d\sigma = -\int_{\tilde{\Sigma}_{2}} \langle W, \nabla \varphi \rangle d\sigma.$$
(8.22)

We would like to set  $\varphi = p$  in (8.22).

Consider two cases:

(i)  $\chi_+ = 1$  a.e. on  $\tilde{\Sigma}_3 \Rightarrow p = 1$  in  $\tilde{\Omega} \Rightarrow p \in C^1(\tilde{\Omega})$ .

(ii)  $\chi_{+} = 0$  on a subset of  $\tilde{\Sigma}_{3}$  of nonzero measure.

Case (i) corresponds to complete saturation of the constraint on the boundary  $\tilde{\Sigma}_3$ . This cannot occur under Hypothesis 6.1. As for case (ii), we know that under Hypothesis 6.2, there exists a neighborhood N of the curve  $C_3 = \tilde{\Sigma}_2 \cap \tilde{\Sigma}_3$  such that  $p = \chi_+ = 0$  on  $\overline{N} \cap \tilde{\Sigma}_3$  ( $\overline{N}$ , the closure of N). This makes it possible to show that p belongs to  $H^{3/2+\sigma}(V_2)$  in some neighborhood  $V_2$  of  $\tilde{\Sigma}_2$ . By construction, p globally belongs to  $H^{1-\epsilon}(\tilde{\Omega})$  for all  $\epsilon > 0$ . However, in order to substitute p into (8.22) we only need to show that the restriction of p to  $V_2$  belongs to  $H^{3/2+\sigma}(V_2)$  for some neighborhood  $V_2$  of  $\tilde{\Sigma}_3$  such that  $\tilde{\Sigma}_3 \cap V_2 \subset N$ . This neighborhood does not contain the discontinuity of the characteristic function  $\chi_+$ ; so the solution p belongs to  $H^{3/2+\sigma}(V_2)$  in that neighborhood of  $\tilde{\Sigma}_2$ .

Setting  $\varphi = p$  in (8.22):

$$\int_{\tilde{\Sigma}_{2}} \frac{\partial Y'}{\partial \nu_{\mathcal{A}}} p |_{\tilde{\Sigma}_{2}} d\sigma = -\int_{\tilde{\Sigma}_{2}} \left[ \tilde{L}^{2} \left( \frac{\partial y}{\partial \xi_{1}} \frac{\partial p}{\partial \xi_{1}} + \frac{\partial y}{\partial \xi_{2}} \frac{\partial p}{\partial \xi_{2}} \right) + \frac{\partial y}{\partial \zeta} \frac{\partial p}{\partial \zeta} \right] \langle V, \nu \rangle d\sigma.$$
(8.23)

By combining (8.20) and (8.23) we finally obtain

$$\lambda(Y') = B(y, p), \tag{8.24}$$

where B is as defined in (8.4).

**Remark 8.2.** In deriving expression (8.24) we have assumed that y at least belonged to  $H^3(\tilde{\Omega})$ . This meant that  $\tilde{\rho}$  was sufficiently smooth and that the angles  $\theta_1$  and  $\theta_3$  were sufficiently small. However, in the end the bilinear form B(y, p) is continuous and well-defined for y and p in  $H^{3/2+\sigma}(V)$ ,  $\sigma > 0$ , for some neighborhood V of  $\tilde{\Sigma}_2$ . In view of the discussion at the end of section 3.4, this means that we can relax conditions on angles  $\theta_1$  and  $\theta_3$  and keep the following ones:

 $0 < \theta_1 < \pi, \qquad 0 < \theta_3 < \pi.$ 

### 8.3. Computation of $\lambda(y_L)$

Recall that the smoothness of  $y_L$  (Remark 7.1) is the same as the smoothness of y provided that the boundary  $\tilde{\Sigma}_2$  is sufficiently smooth and that the angles  $\theta_1$  and  $\theta_3$  are small enough. Assume that y belongs to  $H^2(\tilde{\Omega})$ . Thus  $y_L$  belongs to  $H^2(\tilde{\Omega}) \cap H_0(\tilde{\Omega})$  and we can set  $\psi = y_L$  in (8.11):

$$\lambda(y_L) = (Ay_L, p)_{L^2(\tilde{\Omega})} + \left(\frac{\partial y_L}{\partial \nu_A}, p\right)_{L^2(\tilde{\Sigma}_1)} + \left(\frac{\partial y_L}{\partial \nu_A}, p\right)_{L^2(\tilde{\Sigma}_2)}.$$
(8.25)

In view of identities (7.3)–(7.4) (resp. (7.3a)–(7.4a)) for  $y_L$  the above expression

reduces to

$$\lambda(y_L) = 2\tilde{L} \int_{\tilde{\Omega}} \left( \frac{\partial^2 y}{\partial \xi_1^2} + \frac{\partial^2 y}{\partial \xi_2^2} \right) p \, d\xi \, d\zeta - 2\tilde{L} \int_{\tilde{\Sigma}_2} \left( \nu_1 \frac{\partial y}{\partial \xi_1} + \nu_2 \frac{\partial y}{\partial \xi_2} \right) p \, d\sigma \quad (8.26)$$

$$\lambda(y_L) = + \frac{2}{\tilde{L}} \left[ -\int_{\tilde{\Omega}} \frac{\partial^2 y}{\partial \zeta^2} p \, d\xi \, d\zeta + \int_{\tilde{\Sigma}_2} \nu_{\zeta} \frac{\partial y}{\partial \zeta} p \, d\sigma \right].$$
(8.26a)

**Remark 8.3.** If somehow we "regularize" the characteristic function  $\chi_+$  so that p belongs to  $H^1(\hat{\Omega})$ , the right hand side of expressions (8.26) and (8.26a) can be integrated by parts to yield

$$\lambda(y_L) = -2\tilde{L} \int_{\tilde{\Omega}} \left( \frac{\partial y}{\partial \xi_1} \frac{\partial p}{\partial \xi_1} + \frac{\partial y}{\partial \xi_2} \frac{\partial p}{\partial \xi_2} \right) d\xi d\zeta = b(y, p)$$
(8.27)

(cf. notation (8.3)),

$$\lambda(y_L) = \frac{2}{\tilde{L}} \int_{\tilde{\Omega}} \frac{\partial y}{\partial \zeta} \frac{\partial p}{\partial \zeta} d\xi d\zeta.$$
(8.27a)

**Remark 8.4.** Expression (8.26) only requires that y belongs to  $H^2(\tilde{\Omega}) \cap H_0(\tilde{\Omega})$ , that the boundary  $\tilde{\Sigma}_2$  be smooth (say  $\tilde{\rho}$  belongs to  $C^1(0, 1)$ ) and that the angles  $\theta_1$  and  $\theta_3$  be such that

$$0 < \theta_1 < \pi, \qquad 0 < \theta_3 \leq \frac{\pi}{2}.$$

When p is "regularized" in such a way that it belongs to  $H^1(\tilde{\Omega})$ , the last requirement on  $\theta_3$  can be relaxed to the one

$$0 < \theta_3 < \pi.$$

# 8.4. Final Expressions for Gradients in Cylindrical Coordinates

Let Assumption 6.1 and 8.1 be verified. In cylindrical coordinates, expressions (8.2), (8.3) and (8.4) become

$$a(\psi,\varphi) = \int_0^1 d\zeta \int_0^{\bar{\rho}(\zeta)} 2\pi\rho \, d\rho \left[ \tilde{L}^2 \frac{\partial\psi}{\partial\rho} \frac{\partial\varphi}{\partial\rho} + \frac{\partial\psi}{\partial\zeta} \frac{\partial\varphi}{\partial\zeta} \right]$$
(8.28)

$$b(\psi,\varphi) = -2\tilde{L} \int_0^1 d\zeta \int_0^{\tilde{\rho}(\zeta)} 2\pi\rho \, d\rho \, \frac{\partial\psi}{\partial\rho} \, \frac{\partial\varphi}{\partial\rho} \tag{8.29}$$

$$B(\psi,\varphi) = -\int_0^1 \left[ \tilde{L}^2 \frac{\partial \psi}{\partial \rho} \frac{\partial \varphi}{\partial \rho} + \frac{\partial \psi}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] \omega(\tilde{\rho}(\zeta),\zeta) 2\pi \tilde{\rho}(\zeta) d\zeta.$$
(8.30)

Recall that

$$\lambda(Y') = B(y, p), \tag{8.31}$$

where p is the solution of the boundary-value problem (8.5) and y is the solution of the variational problem: to find y in  $H_0(\tilde{\Omega})$  such that

$$a(y,\varphi) = l(\varphi), \quad \forall \varphi \in H_0(\tilde{\Omega}),$$

$$(8.32)$$

where

$$l(\varphi) = 2\pi \int_0^1 \varphi(0,\rho) \rho \, d\rho.$$
 (8.33)

Expression (6.24) for the gradient of the penalized cost with respect to the shape of the volume  $\tilde{\Omega}$  can now be rewritten in terms of p:

$$\frac{1}{2\pi} dJ_{\varepsilon}(\tilde{L}, \tilde{\rho}; V) = \int_{0}^{1} \left\{ \tilde{L} - \frac{1}{\varepsilon} \left[ \tilde{L}^{2} \frac{\partial y}{\partial \rho} \frac{\partial p}{\partial \rho} + \frac{\partial y}{\partial \zeta} \frac{\partial p}{\partial \zeta} \right]_{\tilde{\Sigma}_{2}} \right\} \tilde{\rho}(\zeta) \omega(\tilde{\rho}(\zeta), \zeta) d\zeta$$
(8.34)

Remark 8.5. Notice that the boundary conditions

$$\frac{\partial p}{\partial \boldsymbol{\nu}_{A}}\Big|_{\boldsymbol{\Sigma}_{2}} = \frac{\partial y}{\partial \boldsymbol{\nu}_{A}}\Big|_{\boldsymbol{\Sigma}_{2}} = 0$$

can be used to obtain

$$B(y,p) = -\int_0^1 \left[1 + \left(\frac{1}{\tilde{L}}\frac{d\tilde{\rho}(\zeta)}{d\zeta}\right)^2\right] \frac{\partial y}{\partial \zeta}\Big|_{\tilde{\Sigma}_2} \frac{\partial p}{\partial \zeta}\Big|_{\tilde{\Sigma}_2} \omega(\tilde{\rho}(\zeta),\zeta) 2\pi\tilde{\rho}(\zeta) d\zeta.$$

This alternate expression can also be used in (8.34).

Similarly recall that

$$\lambda(y_L) = -\frac{4\pi}{\tilde{L}} \left[ \int_0^1 d\zeta \int_0^{\tilde{\rho}(\zeta)} \rho \, d\rho \frac{\partial^2 y}{\partial \zeta^2} p + \int_0^1 d\zeta \tilde{\rho}(\zeta) \frac{d\tilde{\rho}(\zeta)}{d\zeta} \left. \frac{\partial y}{\partial \zeta} \right|_{\tilde{\Sigma}_2} p|_{\tilde{\Sigma}_2} \right] (8.35)$$

(the latter expression for  $\lambda(y_L)$  is preferred to the equivalent one

$$\lambda(y_L) = 4\pi \tilde{L} \left[ \int_0^1 d\zeta \int_0^{\tilde{\rho}(\zeta)} d\rho \frac{\partial}{\partial \rho} \left( \rho \frac{\partial y}{\partial \rho} \right) p - \int_0^1 d\zeta \tilde{\rho}(\zeta) \frac{\partial y}{\partial \rho} |_{\tilde{\Sigma}_2} p|_{\tilde{\Sigma}_2} \right]. \quad (8.35a)$$

Expression (7.8) for the derivative of the penalized cost with respect to the parameter  $\tilde{L}$  can also be written in terms of p:

$$\frac{1}{2\pi} \frac{dJ}{d\tilde{L}} \epsilon(\tilde{L}, \tilde{\rho}) = \int_0^1 \left\{ \frac{1}{2} \tilde{\rho}(\zeta)^2 + \left[ \frac{2}{\epsilon \tilde{L}} \tilde{\rho}(\zeta) \frac{d\tilde{\rho}}{d\zeta}(\zeta) \frac{\partial y}{\partial \zeta} |_{\tilde{\Sigma}_2} p |_{\tilde{\Sigma}_2} \right] + \int_0^{\tilde{\rho}(\zeta)} \rho \frac{\partial^2 y}{\partial \zeta^2} p \, d\rho d\zeta \right\} d\zeta.$$
(8.36)

If p is "regularized" to have it in  $H^1(\tilde{\Omega})$ . Then p is the solution of the variational problem: to find p in  $H_{\chi}(\tilde{\Omega})$  such that

$$a(\psi, p) = 0, \quad \forall \psi \in H_0(\tilde{\Omega})$$

$$(8.37)$$

where

$$H_{\chi}(\tilde{\Omega}) = \left\{ \psi \in H^{1}(\tilde{\Omega}) : \psi|_{\tilde{\Sigma}_{3}} = \chi^{R}_{+} \right\}$$
(8.38)

and  $\chi^{R}_{+}$  is the "regularized"  $\chi_{+}$ . Moreover

$$\lambda(y_L) = b(y, p) = \frac{2}{\tilde{L}} \int_0^1 d\zeta \int_0^{\tilde{\rho}(\zeta)} 2\pi\rho \, d\rho \frac{\partial y}{\partial \zeta} \frac{\partial p}{\partial \zeta}$$
$$= -2\tilde{L} \int_0^1 d\zeta \int_0^{\tilde{\rho}(\zeta)} 2\pi\rho \, d\rho \frac{\partial y}{\partial \rho} \frac{\partial p}{\partial \rho}$$
(8.39)

and

$$\frac{1}{2\pi} \frac{dJ_{\epsilon}}{d\tilde{L}} (\tilde{L}, \tilde{\rho}) = \int_{0}^{1} \left[ \frac{1}{2} \tilde{\rho}(\zeta)^{2} - \frac{2}{\epsilon \tilde{L}} \int_{0}^{\tilde{\rho}(\zeta)} \rho \frac{\partial y}{\partial \zeta} \frac{\partial p}{\partial \zeta} d\rho \right] d\zeta.$$
(8.40)

# 9. Solution of the Penalized Problem

In this section we describe the method which will be used to compute and approximate the solution to the minimization problem

$$Inf\{J_{\varepsilon}(\tilde{L},\tilde{\rho}):\tilde{L} \ge 0, \tilde{\rho} \text{ satisfying (9.2)}\},$$
(9.1)

where

$$\tilde{\rho}(\zeta) \ge 0, \qquad 0 \le \zeta \le 1, \qquad \tilde{\rho}(0) = 1, \qquad \tilde{\rho}(1) \ge \sqrt{1/q}.$$
(9.2)

Problem (9.1) is first rewritten in the following equivalent form

$$\inf\{\inf\{J_{\epsilon}(\tilde{L},\tilde{\rho}):\tilde{\rho} \text{ satisfying } (9.2)\}:\tilde{L} \ge 0\}.$$
(9.3)

Our original problem (9.1) is then divided into the following two subproblems:

I) given  $\tilde{L}$ , find  $\tilde{\rho}_{\varepsilon}(\tilde{L})$  satisfying (9.2) such that

$$\forall \tilde{\rho} \text{ satisfying (9.2)}, \qquad J_{\varepsilon} (\tilde{L}, \tilde{\rho}_{\varepsilon} (\tilde{L})) \leqslant J_{\varepsilon} (\tilde{L}, \tilde{\rho}), \qquad (9.4)$$

II) Find  $\tilde{L}_{\epsilon} \ge 0$  such that

$$\forall \tilde{L} \ge 0, \qquad J_{\varepsilon} \big( \tilde{L}_{\varepsilon}, \tilde{\rho}_{\varepsilon} \big( \tilde{L}_{\varepsilon} \big) \big) \leqslant J_{\varepsilon} \big( \tilde{L}, \tilde{\rho}_{\varepsilon} \big( \tilde{L} \big) \big). \tag{9.5}$$

Problem II will be solved by a one-dimensional search. So we only concentrate on problem I) for a fixed  $\tilde{L}$ . We have seen that expression (8.34) gives the "directional derivative" of  $J_e(\tilde{\Omega})$  with respect to the "speed" or "velocity field"

$$V(\rho,\zeta) = \begin{pmatrix} \omega(\rho,\zeta) \\ 0 \end{pmatrix}.$$
(9.6)

It is of the form (cf. (8.34))

$$j(v) = \int_0^1 f^0(\zeta) v(\zeta) \, d\zeta,$$
(9.7)

where

$$v(\zeta) = \omega(\tilde{\rho}(\zeta), \zeta)$$
(9.8)

anđ

$$f^{0}(\zeta) = \left\{ \tilde{L} - \frac{1}{\varepsilon} \left[ \tilde{L}^{2} \frac{\partial y}{\partial \rho} \frac{\partial p}{\partial \rho} + \frac{\partial y}{\partial \zeta} \frac{\partial p}{\partial \zeta} \right] \right\} 2 \pi \tilde{\rho}(\zeta).$$
(9.9)

The computation of  $f^0$  involves the computation of the state y and the adjoint state p.

## 9.1. Approximation of y and p

We use a finite element method to approximate the solution y of problem (8.32) and the solution p of problem (8.5). Since both solutions are defined on the same domain  $\tilde{\Omega}$ , we shall use the same triangulation of the domain D which generates the domain  $\tilde{\Omega}$  by revolution about the  $\zeta$ -axis.

The triangulation of the domain D will be obtained by deformation of a triangulation of the unit square (cf. Fig. 5)

$$D_0 = \langle (\rho, \zeta) | 0 < \rho < 1, 0 < \zeta < 1 \rangle.$$
(9.10)

Partition the  $\zeta$ -axis into  $M \ge 1$  intervals defined by a sequence  $\{\zeta_j : j = 1, 2, ..., M\}$ 



Fig. 5. Triangulation of the unit square  $D_0$ .

+1 of real numbers such that

$$\zeta_1 = 0, \qquad \zeta_{j+1} > \zeta_j, \qquad j = 1, 2, \dots, M, \qquad \zeta_{M+1} = 1.$$
 (9.11)

Define the parameters

$$k = \max\{\zeta_{i+1} - \zeta_i | j = 1, \dots, M\}.$$
(9.12)

Partition the  $\rho$ -axis into  $N \ge 1$  intervals defined by the sequence  $\{\rho_i = (i-1)/N | i = 1, ..., N+1\}$  of real numbers. Consider the points

$$P_{ij} = (\rho_i, \zeta_j), \quad i = 1, 2, \dots, N+1, \quad j = 1, 2, \dots, M+1.$$
 (9.13)

A triangulation of the domain  $D_0$  is obtained by considering the set of all triangles with vertices

$$(P_{ij}, P_{i+1,j}, P_{i,j+1}), (P_{i,j+1}, P_{i+1,j}, P_{i+1,j+1})$$
 (9.14)

for i = 1, 2, ..., N + 1 and j = 1, 2, ..., M + 1.

The boundary  $\tilde{\Sigma}_2$  of  $\tilde{\Omega}$  and, a fortiori, the boundary  $S_2$  of D, is completely specified by the shape function  $\tilde{\rho}$ . In the finite element approach the boundary  $S_2$ has first to be approximated by a polygonal curve. This means that the function  $\tilde{\rho}$ is approximated by a continuous shape function  $\tilde{\rho}_k$  which is linear on each interval  $[\zeta_j, \zeta_{j+1}], j = 1, 2, ..., M$ . Equivalently,  $\tilde{\rho}_k$  is completely characterized by a sequence  $\{\tilde{\rho}_i | \tilde{\rho}_i > 0, j = 2, ..., M + 1, \tilde{\rho}_1 = 1\}$  of real numbers in the following way

$$\tilde{\rho}_{k}(\zeta) = \frac{\zeta_{j+1} - \zeta}{\zeta_{j+1} - \zeta_{j}} \tilde{\rho}_{j} + \frac{\zeta - \zeta_{j}}{\zeta_{j+1} - \zeta_{j}} \tilde{\rho}_{j+1}, \qquad \zeta \in [\zeta_{j}, \zeta_{j+1}], \tag{9.15}$$

for j = 1, ..., M. This defines a polygonal domain

$$D_k = \left\{ \left(\rho, \zeta\right) | 0 < \zeta < 1, 0 < \rho < \tilde{\rho}_k(\zeta) \right\}$$
(9.16)

which is an approximation of the original domain D.

The triangulation of the domain  $D_k$  is obtained from the triangulation of the unit square  $D_0$  by moving each vertex  $P_{ii}$  in  $D_0$  to the new position (cf. Fig. 6)

$$\tilde{P}_{ij} = \left(\tilde{\rho}_j \rho_i, \zeta_j\right) = \left((i-1)\tilde{\rho}_j / N, \zeta_j\right)$$
(9.17)

in the  $(\rho, \zeta)$ -plane for j = 1, 2, ..., M + 1 and i = 1, 2, ..., N + 1. Notice that the triangulation of  $D_0$  has been chosen in such a way that the resulting triangles in the domain  $D_k$  are not too much distorted. Associate with the triangulation of  $D_k$  the parameter

$$h = \max\{k, k'\}, k' = \max\{\tilde{\rho}_j / N, j = 1, ..., M + 1\}.$$
(9.18)



Fig. 6. Triangulation of the domain  $D_k$ .

Define the finite dimensional subspace  $V_0^h(D_k)$  of  $V_0(D_k)$  as

$$V_0^h(D_k) = \left\{ v_h | v_h \in C^0(\overline{D}_k), v_h \text{ linear on each element,} \right.$$
$$v_h(\rho, 1) = 0, 0 \le \rho \le \tilde{\rho}_k(1) \right\}$$
(9.19)

(the dimension of this space will be M(N+1)). The approximation  $y_h$  of y is defined as the solution of

$$y_{h} \in V_{0}^{h}(D_{k}) \text{ such that for all } v_{h} \text{ in } V_{0}^{h}(D_{k})$$

$$\int_{D_{k}} \left( \tilde{L}^{2} \frac{\partial y_{h}}{\partial \rho} \frac{\partial v_{h}}{\partial \rho} + \frac{\partial y_{h}}{\partial \zeta} \frac{\partial v_{h}}{\partial \zeta} \right) \rho \, d\rho \, d\zeta = \int_{0}^{1} v_{h}(\rho, 0) \rho \, d\rho.$$
(9.20)

By expressing  $y_h$  with respect to a basis of  $V_0^h(D_k)$ , problem (9.20) yields a linear system of algebraic equations and its (unique) solution provides the components of  $y_h$ .

The approximation  $p_h$  of p is obtained by first regularizing the characteristic function  $\chi_+$  in order to put the adjoint variable in the space  $H^1(\tilde{\Omega})$ . Then  $p_h$  is computed in a similar way to  $y_h$  by now using the variational problem (8.37)–(8.38). Choose as the regularized function  $\chi_+^R$  the function  $\chi_+^h$  which is defined by assigning to the points  $\hat{\rho}_i = \tilde{\rho}_{M+1}\rho_i$ , i = 1, ..., N+1, the values

$$\chi_{+}^{h}(\hat{\rho}_{i}) = \begin{cases} 1, & \text{if}\left(\frac{\partial y_{h}}{\partial \zeta}(\hat{\rho}_{i},1)+q\right) < 0\\ 0, & \text{if}\left(\frac{\partial y_{h}}{\partial \zeta}(\hat{\rho}_{i},1)+q\right) \ge 0 \end{cases}$$
(9.21)

and letting  $\chi_{+}^{h}$  be continuous on the whole interval  $[0, \hat{\rho}_{N+1}] = [0, \tilde{\rho}_{M+1}]$  and linear on each interval  $[\hat{\rho}_{i}, \hat{\rho}_{i+1}], i = 1, ..., N$ .

Associate with  $\chi^h_+$  the function space

$$V_{\chi}^{h}(D_{k}) = \left\{ v_{h} \middle| \begin{array}{l} v_{h} \in C^{0}(\overline{D}_{k}) v_{h} \text{ linear on each element} \\ v_{h}(\rho, 1) = \chi_{+}^{h}(\rho), 0 \leq \rho \leq \tilde{\rho}_{k}(1) \end{array} \right\}.$$

Define  $p_h$  as the solution of the variational problem

$$p_{h} \in V_{\chi}^{h}(D_{k}) \text{ such that for all } v_{h} \text{ in } V_{0}^{h}(D_{k})$$

$$\int_{D_{k}} \left( \tilde{L}^{2} \frac{\partial p_{h}}{\partial \rho} \frac{\partial v_{h}}{\partial \rho} + \frac{\partial p_{h}}{\partial \zeta} \frac{\partial v_{h}}{\partial \zeta} \right) \rho \, d\rho \, d\zeta = 0.$$
(9.23)

Problem (9.23) also leads to an algebraic system which can be solved in an analogous way to problem (9.20).

### 9.2. Gradient Computations for the Approximate Problem

In the discretization process we have restricted our attention to domains  $\tilde{\Omega}_h$  which are defined by revolution of the surface  $D_k$  around the  $\zeta$ -axis. The surface  $D_k$  is itself completely specified by the function  $\tilde{\rho}_k$  on [0, 1] which belongs to the subspace

$$P_1^1 = \left\{ \tilde{\rho} \in L^2(0,1) \middle| \begin{array}{l} \tilde{\rho} \text{ is continuous on } [0,1], \tilde{\rho}(0) = 1 \\ \tilde{\rho} \text{ is linear on } [\zeta_j, \zeta_{j+1}], j = 1, \dots, M \end{array} \right\}.$$
(9.24)

We discretize expression (6.24) which gives the Eulerian derivative of the penalized cost for the speed  $V = (\omega, 0)$ :

$$dJ_{\varepsilon}(\tilde{L},\tilde{\rho};V) = 2\pi \left\{ \tilde{L} \int_{0}^{1} \tilde{\rho}(\zeta) \omega(\tilde{\rho}(\zeta),\zeta) d\zeta - \frac{1}{\varepsilon} \int_{0}^{\tilde{\rho}(1)} \frac{\partial Y'}{\partial \zeta}(\rho,1) \chi_{+}(\rho) \rho d\rho \right\}$$
(9.25)

The "discretized speeds"  $V_h = (\omega_h, 0)$  are chosen in such a way that on  $\tilde{\Sigma}_2^h$  (the lateral boundary of  $\tilde{\Omega}_h$ )

$$\omega_h(\tilde{\rho}(\zeta),\zeta) = v(\zeta), \tag{9.26}$$

where v belongs to the subspace

$$P_0^1 = \left\{ v \in L^2(0,1) \middle| \begin{array}{l} v \text{ is continuous on } [0,1], v(0) = 0\\ v \text{ is linear on } [\zeta_j, \zeta_{j+1}], j = 1, \dots, M \end{array} \right\}$$
(9.27)

The variable Y' (solution of (5.16)) is approximated by  $Y'_h$  the solution in  $V_0^h(D_k)$  of the variational problem

$$\int_{D_{k}} \left[ \tilde{L}^{2} \frac{\partial Y_{h}'}{\partial \rho} \frac{\partial \varphi}{\partial \rho} + \frac{\partial Y_{h}'}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] \rho \, d\rho \, d\zeta$$
$$= -\int_{0}^{1} \left[ \tilde{L}^{2} \frac{\partial y}{\partial \rho} \frac{\partial \varphi}{\partial \rho} + \frac{\partial y}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] 2\pi \tilde{\rho}(\zeta) \omega(\tilde{\rho}(\zeta), \zeta) \, d\zeta$$
(9.28)

for all  $\varphi$  in  $V_0^h(D_k)$ . The normal derivative  $\partial Y' / \partial \nu_A$  on  $\tilde{\Sigma}_3$  (which is equal to  $(\partial Y' / \partial \zeta)(\rho, 1)$ ) is approximated in a weak sense by using Green's formula

$$\int_{\tilde{\Sigma}} \frac{\partial Y'}{\partial \nu_{A}} \varphi \, d\sigma = \int_{\tilde{\Omega}} \left[ \tilde{L}^{2} \left( \frac{\partial Y'}{\partial \xi_{1}} \frac{\partial \varphi}{\partial \xi_{1}} + \frac{\partial Y'}{\partial \xi_{2}} \frac{\partial \varphi}{\partial \xi_{2}} \right) + \frac{\partial Y'}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] d\tilde{\Omega} + \int_{\tilde{\Omega}} - A Y' \varphi \, d\Omega \tag{9.29}$$

and the characterization (5.18)-(5.24) of Y'. In the end

$$\int_{\tilde{\Sigma}_{3}} \frac{\partial Y'}{\partial \nu_{\mathcal{A}}} \varphi \, d\sigma = \int_{\tilde{\Omega}} \left[ \tilde{L}^{2} \left( \frac{\partial Y'}{\partial \xi_{1}} \frac{\partial \varphi}{\partial \xi_{1}} + \frac{\partial Y'}{\partial \xi_{2}} \frac{\partial \varphi}{\partial \xi_{2}} \right) + \frac{\partial Y'}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] d\tilde{\Omega} \\ + \int_{\tilde{\Sigma}_{2}} \left[ \tilde{L}^{2} \left( \frac{\partial y}{\partial \xi_{1}} \frac{\partial \varphi}{\partial \xi_{1}} + \frac{\partial y}{\partial \xi_{2}} \frac{\partial \varphi}{\partial \xi_{2}} \right) + \frac{\partial y}{\partial \zeta} \frac{\partial \varphi}{\partial \zeta} \right] \langle V, \nu \rangle \, d\sigma \qquad (9.30)$$

for all  $\varphi$  in  $H^{3/2+\sigma}(\tilde{\Omega})$  such that  $\varphi = 0$  on  $\tilde{\Sigma}_2 \cap \tilde{\Sigma}_3$ . Using the approximation  $Y'_h$  of Y' and speeds  $V_h$ , the normal derivative  $\frac{\partial Y'}{\partial \nu_A}$  is approximated by  $\left(\frac{\partial Y'}{\partial \nu_A}\right)_h$  defined as

$$\begin{split} \int_{(\tilde{\Sigma}_{3})_{h}} \left( \frac{\partial Y'}{\partial \nu_{A}} \right)_{h} \varphi_{h} d\sigma &= \int_{\tilde{\Omega}_{h}} \left[ \tilde{L}^{2} \left( \frac{\partial Y'_{h}}{\partial \xi_{1}} \frac{\partial \varphi_{h}}{\partial \xi_{1}} + \frac{\partial Y'_{h}}{\partial \xi_{2}} \frac{\partial \varphi_{h}}{\partial \xi_{2}} \right) + \frac{\partial Y'_{h}}{\partial \xi} \frac{\partial \varphi_{h}}{\partial \xi} \right] d\tilde{\Omega}_{h} \\ &+ \int_{(\tilde{\Sigma}_{2})_{h}} \left[ \tilde{L}^{2} \left( \frac{\partial y_{h}}{\partial \xi_{1}} \frac{\partial \varphi_{h}}{\partial \xi_{1}} + \frac{\partial y_{h}}{\partial \xi_{2}} \frac{\partial \varphi_{h}}{\partial \xi_{2}} \right) + \frac{\partial y_{h}}{\partial \xi} \frac{\partial \varphi_{h}}{\partial \xi} \right] \\ &\times \langle V_{h}, \nu \rangle d\sigma \end{split}$$
(9.31)

for all  $\varphi_h$  in  $V^h(D_k)$ ,

$$V^{h}(D_{k}) = \left\{ v_{h} | v_{h} \in C^{0}(\overline{D}_{k}), v_{h} \text{ linear on each element} \right\}$$
(9.32)

such that  $\varphi_h = 0$  on  $(\tilde{\Sigma}_2)_h \cap (\tilde{\Sigma}_3)_h$ . Finally the approximation of the Eulerian derivative of the penalized cost for speeds  $V_h$  is given by the expression

$$dJ_{\varepsilon}^{h}(\tilde{L},\tilde{\rho}_{k};V^{h}) = 2\pi \left\{ \tilde{L} \int_{0}^{1} \tilde{\rho}_{k}(\zeta) \omega_{h}(\tilde{\rho}_{k}(\zeta),\zeta) d\zeta - \frac{1}{\varepsilon} \int_{0}^{\tilde{\rho}_{k}(1)} \left(\frac{\partial Y'}{\partial \zeta}\right)_{h}(\rho,1) \chi_{+}(\rho) \rho d\rho.$$
(9.33)

If we assume that  $\chi_+$  is zero on  $\tilde{\Sigma}_3$  in a neighborhood of  $\tilde{\Sigma}_2 \cap \tilde{\Sigma}_3$  and if we regularize  $\chi_+$  to  $\chi_+^h$ , it is easy to verify that expression (9.32) becomes

$$dJ^{h}_{\varepsilon}(\tilde{L},\tilde{\rho}_{k};V^{h}) = \int_{0}^{1} f^{0}_{h}(\zeta) v(\zeta) d\zeta$$
(9.34)

where

$$f_h^0(\zeta) = \left\{ \tilde{L} - \frac{1}{\varepsilon} \left[ \tilde{L}^2 \frac{\partial y_h}{\partial \rho} \frac{\partial p_h}{\partial \rho} + \frac{\partial y_h}{\partial \zeta} \frac{\partial p_h}{\partial \zeta} \right]_{(\tilde{\rho}_k(\zeta),\zeta)} \right\} 2\pi \tilde{\rho}_k(\zeta)$$
(9.35)

(substitute  $\varphi_h = p_h$ , the solution of (9.23), directly into (9.31)).

# 9.3. Choice of the Speed

Once the domain  $D_k$  has been fixed by specifying  $\tilde{\rho}_k$ , the state  $y_h$  and the adjoint state  $p_h$  can be computed. Then the gradient of the penalized cost  $J_{\varepsilon}^h(\tilde{L}, \tilde{\rho}_k)$  can be obtained by computing  $f_h^0(\zeta)$  from equation (9.35).

The approximate gradient in the direction of the speed v is of the form

$$j(v) = \int_0^1 f_h^0(\zeta) v(\zeta) \, d\zeta.$$
(9.36)

Ideally the speed which defines a direction of steepest descent is

$$v(\zeta) = -f_h^0(\zeta) / \|f_h^0\|_2.$$
(9.37)

Unfortunately the speeds v are required to belong to the subspace  $P_0^1$  and the function  $f_h^0$  is generally discontinuous at the points  $\zeta_j$  (the partial derivative of  $y_h$  and  $p_h$  are at best piecewise constant on the boundary  $S_2$ ) and  $f_h^0(0) \neq 0$ . So the best we can do is to project (in the  $L^2$ -sense) the function  $f_h^0$  onto the subspace  $P_0^1$ . More precisely we shall choose the speed

$$v = -\mathcal{P}f_h^0 / \|\mathcal{P}f_h^0\|_2, \tag{9.38}$$

where  $\mathfrak{P}f_h^0$  is the  $L^2$ -projection of  $f_h^0$  onto  $P_0^1$  (cf. Fig. 7).



**Fig. 7.** Gradient  $f^0$  and speed v. First iteration,  $\varepsilon = 1/50$ , L = 2.75.

### 9.4. Method of Descent

Start with an initial shape function  $\tilde{\rho}_k^0$  such that

$$\tilde{\rho}_k^0(\zeta) \ge 0, \qquad 0 \le \zeta \le 1, \qquad \tilde{\rho}_k^0(0) = 1, \qquad \tilde{\rho}_k^0(1) \ge \sqrt{1/q} \,. \tag{9.39}$$

Assume that at step n + 1 the shape function  $\tilde{\rho}_k^n$  has been constructed in such a way that conditions (9.38) are verified with  $\tilde{\rho}^n$  in place of  $\tilde{\rho}^0$ . To construct the new shape function  $\tilde{\rho}_k^{n+1}$ , we construct the new domain  $D_k^n$  from  $\tilde{\rho}_k^n$  and compute the functions  $y_k^n$  and  $p_k^n$  on  $D_k^n$ . We compute the function  $(f_h^0)^n$  and the speed  $v^n$  given by (9.37) with  $(f_h^0)^n$  in place of  $f_h^0$ . Finally we construct the new shape function

$$\tilde{\rho}_{k}^{n+1}(\zeta) = \tilde{\rho}_{k}^{n}(\zeta) + t^{n+1}v^{n}(\zeta)$$
(9.40)

by choosing an appropriate  $t^{n+1} > 0$  which minimizes  $J_{\epsilon}(\tilde{L}, \tilde{\rho}_{k}^{n+1}(\tilde{L}))$  and verifies condition (9.38) with  $\tilde{\rho}_{k}^{n+1}$  in place of  $\tilde{\rho}_{k}^{0}$ .

### 9.5. Numerical Example

Three optimal shapes will be considered here, corresponding to the parameter  $R_0 = 0.57$  cm and the three values 1/5, 1/10 and 1/20 of the flux ratio q.

We shall describe detailed computations for the intermediate case q = 1/10. The discretization parameters are the following:

$$N = 10, \quad M = 14$$
  

$$\{\zeta_j\} = \{0.; 0.05; 0.1; 0.175; 0.25; 0.35; 0.45; 0.55; 0.65; 0.74; 0.82; 0.88; 0.93; 0.97; 1.\}.$$

A typical triangulation of the domain  $D_k$  is shown in Fig. 8 with an example of isotherms giving the temperature field inside the diffuser. For a fixed value of the





Fig. 8. Finite element grid and isotherms.

parameter  $\tilde{L}$ , a solution  $\tilde{\rho}_{\epsilon}(\tilde{L})$  of Problem I is found by the descent method of sec. 6.4. The initial shape  $\tilde{\rho}_{k}^{0}$  is chosen as the linear function

$$\tilde{\rho}_k^0(\zeta) = 1 + 2\zeta/R_0, \quad 0 \le \zeta \le 1, \quad R_0 = 0.57 \text{cm.}$$
 (9.41)

In a first step, the parameter  $\varepsilon$  is set equal to 1/50 in order to let the constraint be active during the iterations (cf. Fig. 9). The penalized cost can be plotted against the length  $L = R_0 \tilde{L}$  as shown in Fig. 10. From this graph it is readily seen that the length L leading to the lightest diffuser is around L = 2.75cm. Details of a typical computation are shown on Fig. 11. We obviously see a monotoneously decreasing penalized cost, a "dual behavior between the volume and the constraint (one increasing and the other decreasing). One of the striking features of the optimization search is that a small difference between two shapes can result in finding or not a "good" direction of descent. Fig. 11 shows the evolution of the shape function during the iteration for a fixed  $L = \tilde{L}R_0 = 2.75$ .





Fig. 11. Results for one optimization of the shape starting from a linear initial shape;  $\varepsilon = 1/50$ .



Fig. 13. Optimal diffusers for different ratios q.

In a second step, the parameter  $\varepsilon$  is reduced from the value 1/50 to 1/250. The initial shape function is chosen as the one found at the end of the iterations of the first step. In so doing we reduce the value of the constraint functional to practically zero by slightly enlarging the diffuser near its output surface  $\tilde{\Sigma}_3$  (see Fig. 12). A one-dimensional search is also performed at this stage to determine the new value of the parameter  $L = \tilde{L}R_0$  (L = 2.775) giving the smallest cost.

The optimal diffusers for the three values of the flux ratio q and  $R_0 = 0.57$  have been drawn on Fig. 13.

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