# Small vibrations of elastic medium deforming in time 

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#### Abstract

Isorropic elastic material is subject to finite strains in such a manner that the elongations in three mutually perpendicular directions are proportional to time. The equations for a small additional motion are constructed and several types of possible vibrations are analyzed. On the basis of the condition of propagation it is demonstrated that three principal directions of propagation exist connected with longitudinal and transversal waves.


Izotropowy materiał sprężysty poddany jest skończonym odksztalceniom w ten sposób, że wydłużenia w trzech wzajemnie prostopadłych kierunkach sa proporcjonalne do czasu. Buduje się równania dla malego dodatkowego ruchu, a nastepnie analizuje kilka możliwych drgań. W oparciu o warunek propagacji pokazuje się, że istnieją trzy główne kierunki propagacji, którym odpowiadają fale poprzeczne lub podłużne.


#### Abstract

Изотропный упругий материал подвержен конечным деформациям таким образом, что удлиннения в трех взаимно перпендикулярных направлениях пропорционально изменяются во времени. Выводятся уравнения, описывающие малое дополнительное движение, которые затем подвергаются анализу с точки зрения различных видов возможных колебаний. На основе условия распространения волн показано, что существуют три основных направления распространения волн, которым соответствуют продольные или поперечные колебания.


Small vibrations of elastic media have been extensively treated in the literature. In particular it is known that a complete coincidence can be established between the theory of vibrations and the condition of propagation of an acoustical wave provided the vibrations are infinitesimal and the finite initial deformation is homogeneous and stationary [1]. Thus, the coincidence also holds true in the linear theory of elasticity where-according to the definition-the initial deformation does not exist. In the present paper, we consider a situation more general than those hitherto dealt with. The finite, initial deformation varies in time; small additional vibrations are now superposed on that deformation and certain particular forms of vibrations are investigated.

## 1. Fundamental motion and additional motion

Let us introduce fixed Cartesian coordinate system. The coordinates of a typical point of the body under consideration in the natural state $B_{R}$ are denoted by $X^{\alpha}, \alpha=1,2,3$. Let us consider the motion $\chi(t)$ given by the relations

$$
\begin{equation*}
x^{1}=\lambda_{1} X^{1}, \quad x^{2}=\lambda_{2} X^{2}, \quad x^{3}=\lambda_{3} X^{3}, \tag{1.1}
\end{equation*}
$$

where $\lambda_{K}$ are certain functions of time $t$ only,

$$
\begin{equation*}
\lambda_{K}=\lambda_{K}(t) \tag{1.2}
\end{equation*}
$$

At the instant $t=0$, the body is in a natural state $B_{R}$, and hence $\lambda_{\mathbf{K}}(0)=1$. Superposition of a rigid translation (but no rotation) upon the motion (1.1) does not influence the subsequent relations of this paper.

Let us pass to determination of the strain gradients $x_{, \alpha,}^{i}$, the strain tensor $B^{i k}$, its invariants $I_{\mathbf{K}}, K=1,2,3$, and the density $\varrho$. All calculations of this section are based on the relations and notation given in [2]. According to (1.1), we have

$$
\begin{gather*}
x_{, \alpha}^{i}=\left[\begin{array}{ccc}
\lambda_{1} & 0 & 0 \\
& \lambda_{2} & 0 \\
& & \lambda_{3}
\end{array}\right],  \tag{1.3}\\
B^{k k}=x_{, \alpha}^{i} \alpha^{k, \alpha}=\left[\begin{array}{ccc}
\lambda_{1}^{2} & 0 & 0 \\
& \lambda_{2}^{2} & 0 \\
& & \lambda_{3}^{2}
\end{array}\right],  \tag{1.4}\\
I_{1}=B_{r}^{r}=\lambda_{1}^{2}+\lambda_{2}^{2}+\lambda_{3}^{2}, \\
I_{2}=\frac{1}{2}\left(I_{1}^{2}-B_{z}^{r} B_{r}^{v}\right)=\lambda_{1}^{2} \lambda_{2}^{2}+\lambda_{2}^{2} \lambda_{3}^{2}+\lambda_{3}^{2} \lambda_{1}^{2},  \tag{1.5}\\
I_{3}=\frac{1}{3} B_{s}^{r} B_{p}^{s} B_{r}^{p}-\frac{1}{2} B_{r}^{r} B_{p}^{s} B_{s}^{p}+\frac{1}{6}\left(B_{r}^{r}\right)^{3}=\lambda_{1}^{2} \lambda_{2}^{2} \lambda_{3}^{2} ; \\
\varrho=\frac{1}{\lambda_{1} \lambda_{2} \lambda_{3}} \varrho_{R}, \tag{1.6}
\end{gather*}
$$

$\varrho_{R}$ being the density in the state $B_{R}$. Raising and lowering of indices is performed by means of the metric tensors of the Cartesian coordinate system previously introduced:

$$
\begin{equation*}
g_{i j}=g_{s}^{i j}=\delta_{i j}, \quad g_{\alpha \beta}=g^{\alpha \beta}=\delta_{\alpha \beta} . \tag{1.7}
\end{equation*}
$$

Further considerations will be confined to homogeneous and isotropic elastic materials. For such materials, there exists the elastic potential $\sigma$ (referred to unit mass) which is a function of the strain invariants $I_{\mathbf{K}}$ only, $\sigma=\sigma\left(I_{\mathbf{K}}\right)$, and the Piola-Kirchhoff stress tensor $T_{R i}{ }^{\alpha}$ is defined by

$$
\begin{equation*}
T_{R i}^{\alpha}=\varrho_{R} \frac{\partial \sigma}{\partial x_{, \alpha}^{i}} \tag{1.8}
\end{equation*}
$$

where $\sigma$ is a function of the gradients $x^{d}$. through the invariants $I_{\mathbf{K}}$. Using the relations following from Eqs. (1.4), (1.5),

$$
\begin{align*}
& \frac{\partial B^{r s}}{\partial x^{i}, \alpha}=\delta_{i}^{r} x^{s \alpha}+\delta_{i}^{s} x^{r \alpha}, \\
& \frac{\partial I_{1}}{\partial B^{r s}}=g^{r s}, \quad \frac{\partial I_{2}}{\partial B^{r s}}=I_{1} g_{r s}-B_{r s},  \tag{1.9}\\
& \frac{\partial I_{3}}{\partial B^{r s}}=B_{r p} B_{s}^{p}-I_{1} B_{r s}+I_{2} g^{r s},
\end{align*}
$$

we obtain

$$
\begin{gather*}
T_{R k}{ }^{\alpha}=2 W_{1} x_{k}{ }^{\alpha}+2 W_{2}\left(I_{1} x_{k},{ }^{\alpha}-B_{k s} x^{s},{ }^{\alpha}\right)+2 W_{3}\left(B_{k p} B^{p q} x_{q}{ }^{\alpha}-I_{1} B_{k p} x^{p},{ }_{\alpha}+I_{2} x_{k}{ }^{\alpha}\right)  \tag{1.10}\\
W_{K}=\frac{\partial W}{\partial I_{K}}, \quad W=\varrho_{R} \sigma, \quad K=1,2,3 .
\end{gather*}
$$

Substituting now Eqs. (1.3), (1.4), (1.5) into Eq. (1.8), we obtain:

$$
\begin{align*}
& T_{R 1}^{1}=2\left[\lambda_{1} W_{1}+\lambda_{1}\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right) W_{2}+\lambda_{1} \lambda_{2}^{2} \lambda_{3}^{2} W_{3}\right], \\
& T_{R 2}^{2}=2\left[\lambda_{2} W_{1}+\lambda_{2}\left(\lambda_{3}^{2}+\lambda_{1}^{2}\right) W_{2}+\lambda_{2} \lambda_{3}^{2} \lambda_{1}^{2} W_{3}\right],  \tag{1.11}\\
& T_{R 3}{ }^{3}=2\left[\lambda_{3} W_{1}+\lambda_{3}\left(\lambda_{1}^{2}+\lambda_{2}^{2}\right) W_{2}+\lambda_{3} \lambda_{1}^{2} \lambda_{2}^{2} W_{3}\right], \\
& T_{R 1}^{2}=T_{R 2}^{1}=T_{R 1}^{3}=T_{R 3}{ }^{1}=T_{R 2}{ }^{3}=T_{R 3}{ }^{2}=0 .
\end{align*}
$$

Since $T_{R i}{ }^{\alpha}$ are independent of the coordinates $X^{\alpha}$, the left-hand side of the equations of motion

$$
\begin{equation*}
T_{R i, \alpha}^{\alpha}=\varrho_{R} \ddot{x}_{i} \tag{1.12}
\end{equation*}
$$

is equal to zero, which yields the conclusion that also the acceleration $\ddot{x}_{i}$ is equal to zero. The motion (1.1) is then possible provided that

$$
\begin{equation*}
\lambda_{1}=1+c_{1} t, \quad \lambda_{2}=1+c_{2} t, \quad \lambda_{3}=1+c_{3} t \tag{1.13}
\end{equation*}
$$

where $c_{1}, c_{2}, c_{3}$ are fixed parameters. In the subsequent considerations it is assumed that Eq. (1.13) has been subsituted in Eq. (1.1).

Let us now pass to the consideration of a perturbed motion $\chi^{*}(t)$, differing only slightly from the motion $\chi(t)$ [cf. Eqs. (1.1) and (1.13)]-that is, the motion

$$
\begin{align*}
& x^{* 1}=\lambda_{1} X^{1}+u^{1}\left(X^{\alpha}, t\right), \\
& x^{* 2}=\lambda_{2} X^{2}+u^{2}\left(X^{\alpha}, t\right),  \tag{1.14}\\
& x^{* 3}=\lambda_{3} X^{3}+u^{3}\left(X^{\alpha}, t\right) .
\end{align*}
$$

The quantity $u^{i}\left(X^{\alpha}, t\right)$ is the displacement of the perturbation. Toupin and Bernstein [3] derived the following equation for $u^{i}$

$$
\begin{equation*}
\left(A_{i}^{\alpha}{ }_{k}^{\beta} u^{k}{ }_{; \beta}^{k}\right)_{; \alpha}=\varrho_{R} \ddot{u}_{i} \tag{1.15}
\end{equation*}
$$

where the functions

$$
\begin{equation*}
A_{i}^{\alpha}{ }_{k}^{\alpha \beta}=\varrho_{R} \frac{\partial^{2} \sigma}{\partial x^{i}{ }_{, \alpha} \partial x^{k}, \beta}, \tag{1.16}
\end{equation*}
$$

are calculated for the fundamental motion $\chi(t)$. In the Cartesian coordinate system introduced here and under a consistent application of the independent variables $X^{\alpha}$ (and not $x^{t}$ ), all differentiations (1.15) may be replaced by partial differentiations with respect to $X^{\alpha}$ and $t$.

In order to obtain Eq. (1.16) in an explicit form in the case of isotropic materials, the necessary differentiations should be performed; it should be born in mind that $\sigma$ depends on $x^{i},{ }_{\alpha}$ through the invariants $I_{\mathbf{K}}$. Applying Eqs. (1.9) once again, we obtain

$$
\begin{equation*}
A_{k}{ }^{\alpha}{ }^{\beta}=2 W_{1} g_{k m} g^{\alpha \beta}+2 W_{2}\left[2 x_{k},{ }^{\alpha} x_{m},{ }^{\beta}-g_{k m} x_{r},{ }^{\alpha} x^{r},{ }^{\beta}-x_{k}{ }^{\beta} x_{m}{ }^{\alpha}+\left(I_{1} g_{k m}-B_{k m}\right) g^{\alpha \beta}\right]+ \tag{1.17}
\end{equation*}
$$

$$
\begin{aligned}
& +2 W_{3}\left[\left(g_{k r} B_{s q} x^{q},{ }^{\alpha}+B_{k r} x_{s,}{ }^{\alpha}-g_{r s} B_{k q} x^{q \alpha} .-I_{1} g_{k r} x_{s,}{ }^{\alpha}+I_{1} g_{r s} x_{k,}{ }^{\alpha}-B_{r s} x_{k}{ }^{\alpha}{ }^{\alpha}\right) \times\right. \\
& \left.\times\left(g_{m}^{r} x^{s \beta}+g_{m}^{s} x^{r}{ }^{\beta}\right)+\left(B_{k p} B_{m}^{p}-I_{1} B_{k m}+I_{2} g_{k m}\right) g^{\alpha \beta}\right] \\
& +4\left\{W_{11} x_{k}{ }^{\alpha} x_{m}{ }^{\beta}+W_{22}\left(I_{1} x_{k}{ }^{\alpha}-B_{k}{ }^{p} x_{p}{ }^{\alpha}{ }^{\alpha}\right)\left(I_{1} x_{m}{ }^{\beta}-B_{m}{ }^{r} x_{r}{ }^{\beta}\right)\right. \\
& +W_{33}\left(B_{k p} B^{p q} x_{q},{ }^{\alpha}-I_{1} B_{k}{ }^{q} x_{q}{ }^{\alpha}{ }^{\alpha}+I_{2} x_{k}{ }^{\alpha}{ }^{\alpha}\right)\left(B_{m r} B^{r s} x_{s,}{ }^{\beta}-I_{1} B_{m}{ }^{r} x_{r},{ }^{\beta}+I_{2} x_{m}{ }^{\beta}\right) \\
& +W_{23}\left[\left(I_{1} x_{k}{ }^{\alpha}{ }^{\alpha}-B_{k}{ }^{p} x_{p,}{ }^{\alpha}\right)\left(B_{m r} B^{r s} x_{s,}{ }^{\beta}-I_{1} B_{m}{ }^{s} x_{s}{ }^{\beta}+I_{2} x_{m,}{ }^{\alpha}\right)\right. \\
& \left.+\left(B_{k p} B^{p q} x_{q,}{ }^{\alpha}-I_{1} B_{k}{ }^{q} x_{q}{ }^{\alpha}{ }^{\alpha}+I_{2} x_{k}{ }^{\alpha}{ }^{\alpha}\right)\left(I_{1} x_{m}{ }^{\beta}-B_{m r} x^{r,}{ }^{\beta}\right)\right] \\
& +W_{31}\left[\left(B_{k p} B^{p q} x_{q,}{ }^{\alpha}-I_{1} B_{k}{ }^{q} x_{q}{ }^{\alpha}{ }^{\alpha}+I_{2} x_{k,}{ }^{\alpha}\right) x_{m,}{ }^{\beta}+x_{k}{ }^{\alpha}{ }^{\alpha}\left(B_{m r} B^{r s} x_{s,}{ }^{\beta}-I_{1} B_{m}{ }^{r} x_{r}{ }^{\beta}+I_{2} x_{m}{ }^{\beta}\right)\right. \\
& +W_{12}\left[\left(I_{1} x_{k},{ }^{\alpha}-B_{k p} x^{p},{ }^{\alpha}\right) x_{m}{ }^{\beta}+x_{k}{ }^{\alpha}\left(I_{1} x_{m}{ }^{\beta}-B_{m r} x^{\sim}, \beta\right)\right], \quad W_{k L}=\frac{\partial^{2} W}{\partial I_{k}, \partial I_{L}} .
\end{aligned}
$$

The functions $A_{i k}$ are symmetric neither in Latin nor Greek indices. Due to Eq. (1.16), however, the symmetry $A_{i k}=A_{k i}$ occurs.

Substituting Eqs. (1.3), (1.4), (1.5) into (1.17), we finally obtain:

$$
\begin{aligned}
& A_{1}{ }^{1}{ }_{1}{ }^{1}=4 \lambda_{1}^{2} W_{11}+4 \lambda_{1}^{2}\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right)^{2} W_{22}+8 \lambda_{1}^{2} \lambda_{2}^{4} \lambda_{3}^{4} W_{33}+8 \lambda_{1}^{2} \lambda_{2}^{2} \lambda_{3}^{2}\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right) W_{23} \\
& +8 \lambda_{1}^{2} \lambda_{2}^{2} \lambda_{3}^{2} W_{31}+4 \lambda_{1}^{2}\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right) W_{12}+2 W_{1}+2\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right) W_{2}+2 \lambda_{2}^{2} \lambda_{3}^{2} W_{3}, \\
& A_{1}{ }^{2}{ }_{1}{ }^{2}=2\left(W_{1}+\lambda_{3}^{2} W_{2}\right), \\
& A_{1}{ }^{3}{ }_{1}{ }^{3}=2\left(W_{1}+\lambda_{2}^{2} W_{2}\right), \\
& A_{1}{ }_{2}{ }^{2}=4 \lambda_{1} \lambda_{2} W_{11}+4 \lambda_{1} \lambda_{2}\left(\lambda_{1}^{2}+\lambda_{3}^{2}\right)\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right) W_{22}+4 \lambda_{1}^{3} \lambda_{2}^{3} \lambda_{3}^{4} W_{33} \\
& +4 \lambda_{1} \lambda_{2} \lambda_{3}^{2}\left(\lambda_{1}^{2} \lambda_{3}^{2}+\lambda_{2}^{2} \lambda_{3}^{2}+2 \lambda_{1}^{2} \lambda_{2}^{2}\right) W_{23}+4 \lambda_{1} \lambda_{2} \lambda_{3}^{2}\left(\lambda_{1}^{2}+\lambda_{2}^{2}\right) W_{31} \\
& +4 \lambda_{1} \lambda_{2}\left(\lambda_{1}^{2}+\lambda_{2}^{2}+2 \lambda_{3}^{2}\right) W_{12}+4 \lambda_{1} \lambda_{2} W_{2}+4 \lambda_{1} \lambda_{2} \lambda_{3}^{2} W_{3}, \\
& A_{1}{ }^{1}{ }_{3}{ }^{3}=4 \lambda_{1} \lambda_{3} W_{11}+4 \lambda_{1} \lambda_{3}\left(\lambda_{1}^{2}+\lambda_{2}^{2}\right)\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right) W_{22}+4 \lambda_{1}^{3} \lambda_{2}^{4} \lambda_{3}^{3} W_{33} \\
& +4 \lambda_{1} \lambda_{2}^{2} \lambda_{3}\left(\lambda_{1}^{2} \lambda_{2}^{2}+\lambda_{2}^{2} \lambda_{3}^{2}+2 \lambda_{3}^{2} \lambda_{1}^{2}\right) W_{23}+4 \lambda_{1} \lambda_{2}^{2} \lambda_{3}\left(\lambda_{1}^{2}+\lambda_{3}^{2}\right) W_{31}, \\
& +4 \lambda_{1} \lambda_{3}\left(\lambda_{1}^{2}+\lambda_{3}^{2}+2 \lambda_{2}^{2}\right) W_{12}+4 \lambda_{1} \lambda_{3} W_{2}+4 \lambda_{1} \lambda_{2}^{2} \lambda_{3} W_{3}, \\
& A_{1}{ }^{2}{ }_{2}{ }^{1}=-2 \lambda_{1} \lambda_{2} W_{2}-2 \lambda_{1} \lambda_{2} \lambda_{3}^{2} W_{3} \text {, } \\
& A_{1}{ }^{3}{ }_{3}{ }^{1}=-2 \lambda_{1} \lambda_{3} W_{2}-2 \lambda_{1} \lambda_{3} \lambda_{2}^{2} W_{3} .
\end{aligned}
$$

All the remaining functions $A_{1}{ }_{k}{ }_{k}^{\beta}$ vanish. The functions $A_{2}{ }_{2}{ }_{k}{ }^{\beta}$ and $A_{3}{ }_{k}{ }_{k}{ }^{\beta}$ may be obtained from those given by cyclic change of indices at the elongations $\lambda_{\mathrm{K}}$.

It is seen that the tensor $A_{i}{ }_{k}^{\alpha}{ }_{k}^{\beta}$ is independent of $\left(x^{k}, X^{\alpha}\right)$ and is the function of the only variable $t$. Owing to that property, the tensor may be taken out of the parantheses in Eq. (1.15). The differentiation indicated in that formula may be reduced-in the Cartesian coordinate system-to partial differentiation (with fixed $X^{\alpha}$ ); hence, we obtain:

$$
\begin{equation*}
A_{i}^{\alpha} k^{\beta} \frac{\partial^{2} u^{k}}{\partial X^{\alpha} \partial X^{\beta}}=\varrho_{R} \frac{\partial^{2} u_{i}}{\partial t^{2}} . \tag{1.19}
\end{equation*}
$$

This equation is the equation sought for, describing the small motion superposed on the fundamental motion (1.1).

## 2. Derivation in convective coordinates

The derivation of Eqs. (1.19) presented above in the case of small additional motion is complete. Since, however, that in solving certain particular problems (e.g., stability), convective coordinates are extensively used in the literature, let us employ these coordinates in the present paper. Such an approach has the advantage that the corresponding formulae are - particularly in the case of isotropic materials-well known from the relevant literature (cf., e.g., [4]).

In addition to the Cartesian reference frames $x^{i}$ and $X^{\alpha}$, let us now introduce the convective coordinate system $\theta^{i}$ which coincides in $B$ with the system $X^{\alpha}$

$$
\begin{equation*}
\theta^{1}=X^{1}, \quad \theta^{2}=X^{2}, \quad \theta^{3}=X^{3} \tag{2.1}
\end{equation*}
$$

The time-dependent metric tensor corresponding to $\theta^{i}$ is denoted by $G^{i j}, G_{i j}$, by contrast with fixed $g_{i j}, g^{i j}, g_{\alpha \beta}, g^{\alpha \beta}$.

Calculations of the invariants $I_{K}$ lead to relations already given (1.5). Passing now to the stress tensor, we obtain

$$
\begin{equation*}
\tau^{11}=\Psi_{1}^{\prime}+\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right) \Psi_{2}+\frac{1}{\lambda_{1}^{2}} \Psi_{3}, \quad \tau^{12}=0 \tag{2.2}
\end{equation*}
$$

with the notations

$$
\begin{equation*}
\Psi_{1}=\frac{2}{\sqrt{I_{3}}} \frac{\partial W}{\partial I_{2}} . \quad \Psi_{2}=\frac{2}{\sqrt{I_{3}}} \frac{\partial W}{\partial I_{2}}, \quad \Psi_{3}=2 \sqrt{I_{3}} \frac{\partial W}{\partial I_{3}}, \tag{2.3}
\end{equation*}
$$

in [3] the corresponding notations were $\Phi, \Psi, p), W$ denoting the elastic energy referred to unit volume in $B_{R}$. This function is identical with the function $W=\varrho_{R} \sigma$ introduced in the preceding section. The remaining components of the stress tensor are obtained from those defined by Eq. (1.11) by means of cyclic interchange of indices.

On the basis of relations given in [3], let us consider the motion $\mathbf{R}^{*}(t)=\mathbf{R}(t)+\mathbf{w}(t)$, $\mathbf{w}$ being small in comparison with $\mathbf{R}$. The quantities appearing in Eqs. (2.2)-(2.3) are now subject to certain increments. Their linear components are denoted by the same kernel letter as those connected with the motion $\mathbf{R}(t)$ and marked by a prime. With the notations

$$
\begin{equation*}
\mathbf{W}=u \mathbf{G}^{1}+v \mathbf{G}^{2}+w \mathbf{G}^{3} \tag{2.4}
\end{equation*}
$$

the physical components of $w$ are, according to (2.2), the quantities $u / \lambda_{1}, v / \lambda_{2}, w / \lambda_{3}$. We now have

$$
\left.\begin{array}{rl}
\tau^{\prime 11}= & \frac{2}{\lambda_{1} \lambda_{2} \lambda_{3}}\left[2 W_{11}\right.
\end{array}+2 W_{22}\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right)^{2}+2 W_{33} \lambda_{2}^{4} \lambda_{3}^{4}+4 W_{23} \lambda_{2}^{2} \lambda_{3}^{2}\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right) ~ 子 W_{31} \lambda_{2}^{2} \lambda_{3}^{2}+4 W_{12}\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right)-W_{1} \frac{1}{\lambda_{1}^{2}}-W_{2} \frac{\lambda_{2}^{2}+\lambda_{3}^{2}}{\lambda_{1}^{2}}-W_{3} \frac{\lambda_{2}^{2} \lambda_{3}^{2}}{\lambda_{1}^{2}}\right] u_{X} .
$$

$$
\begin{align*}
& \begin{aligned}
&\left.-W_{1} \frac{1}{\lambda_{2}^{2}}+W_{2} \frac{\lambda_{2}^{2}-\lambda_{3}^{2}}{\lambda_{2}^{2}}+W_{3} \lambda_{3}^{2}\right] v_{Y}+\frac{2}{\lambda_{1} \lambda_{2} \lambda_{3}}\left[2 W_{11}+2 W_{22}\left(\lambda_{1}^{2}+\lambda_{2}^{2}\right)\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right)\right. \\
&+2 W_{33} \lambda_{1}^{2} \lambda_{2}^{4} \lambda_{3}^{2}+2 W_{23} \lambda_{2}^{2}\left(\lambda_{1}^{2} \lambda_{2}^{2}+\lambda_{2}^{2} \lambda_{3}^{2}+2 \lambda_{1}^{2} \lambda_{3}^{2}\right)+2 W_{31} \lambda_{2}^{2}\left(\lambda_{1}^{2}+\lambda_{3}^{2}\right) \\
&\left.+2 W_{12}\left(\lambda_{1}^{2}+\lambda_{3}^{2}+2 \lambda_{2}^{2}\right)-W_{1} \frac{1}{\lambda_{3}^{2}}+W_{2} \frac{\lambda_{3}^{2}-\lambda_{2}^{2}}{\lambda_{3}^{2}}+W_{3} \lambda_{2}^{2}\right] W_{\mathrm{Z}},
\end{aligned}  \tag{2.5}\\
& \tau^{\prime \prime 12=}=-\frac{2}{\lambda_{1} \lambda_{2} \lambda_{3}}\left(W_{2}+\lambda_{3}^{2} W_{3}\right)\left(u_{Y}+v_{X}\right), \\
& \tau^{\prime 13}=-\frac{2}{\lambda_{1} \lambda_{2} \lambda_{3}}\left(W_{2}+\lambda_{2}^{2} W_{3}\right)\left(u_{z}+w_{X}\right),
\end{align*}
$$

where

$$
\begin{equation*}
W_{K L}=\frac{\partial^{2} W}{\partial I_{K} \partial I_{L}} . \tag{2.6}
\end{equation*}
$$

The remaining increments of the stress tensor are obtained from those given above by cyclic permutation of indices of all quantities except $W_{K L}$.

To the equations of motion constructed by means of convective coordinates, there enter, in addition to the increments already mentioned, the increments of the Christoffel symbols $\Gamma_{j k}^{\prime i}$ and of the acceleration $\mathbf{a}^{\prime}$. In order to determine $\mathbf{a}^{\prime}$, let us differentiate $\mathbf{R}^{*}(t)$ twice with respect to time, which finally yields

$$
\begin{equation*}
a^{\prime 1}=\frac{1}{\lambda_{1}} \frac{D^{2}}{D t^{2}}\left(\frac{u}{\lambda_{1}}\right), \quad a^{\prime 2}=\frac{1}{\lambda_{2}} \frac{D^{2}}{D t^{2}}\left(\frac{v}{\lambda^{2}}\right), \quad a^{\prime 3}=\frac{1}{\lambda_{3}} \frac{D^{2}}{D t^{2}}\left(\frac{w}{\lambda_{3}}\right) . \tag{2.7}
\end{equation*}
$$

Here $D / D t$ denotes the differentiation with respect to $t$ at fixed values of $\theta^{i}$.
Using the formulae derived in [2], the increments of $\Gamma_{j k}^{\prime i}$ are determined (Christoffel symbols $\Gamma_{j k}^{i}$ vanish);

$$
\begin{array}{lll}
\Gamma_{11}^{\prime 1}=\frac{1}{\lambda_{1}^{2}} u_{X X}, & \Gamma_{22}^{\prime 1}=\frac{1}{\lambda_{1}^{2}} u_{Y Y}, & \Gamma_{33}^{\prime 1}=\frac{1}{\lambda_{1}^{2}} u_{Z Z}  \tag{2.8}\\
\Gamma_{23}^{\prime 1}=\frac{1}{\lambda_{1}^{2}} u_{Y Z}, & \Gamma_{31}^{\prime \prime}=\frac{1}{\lambda_{1}^{2}} u_{Z X}, & \Gamma_{12}^{\prime 1}=\frac{1}{\lambda_{1}^{2}} u_{x Y}
\end{array}
$$

The remaining increments $\Gamma_{j k}^{\prime i}$ are obtained from (2.14) by cyclic permutation of indices and $(u, v, w)$.

Inserting now Eqs. (2.7) and (2.8) in the equations of motion we finally obtain:

$$
\begin{align*}
& {\left[2 W_{11}+2 W_{22}\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right)^{2}+2 W_{33} \lambda_{2}^{4} \lambda_{3}^{4}+4 W_{23} \lambda_{2}^{2} \lambda_{3}^{2}\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right)+4 W_{31} \lambda_{2}^{2} \lambda_{3}^{2}\right.}  \tag{2.9}\\
& \left.+4 W_{12}\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right)+\frac{1}{\lambda_{1}^{2}} W_{1}+\frac{\lambda_{2}^{2}+\lambda_{3}^{2}}{\lambda_{1}^{2}} W_{2}+\frac{\lambda_{2}^{2}+\lambda_{3}^{2}}{\lambda_{1}^{2}} W_{3}\right] u_{X X}+\left(\frac{1}{\lambda_{1}^{2}} W_{1}+\frac{\lambda_{3}^{2}}{\lambda_{1}^{2}} W_{2}\right) u_{Y Y} \\
& \\
& +\left(\frac{1}{\lambda_{1}^{2}} W_{1}+\frac{\lambda_{2}^{2}}{\lambda_{1}^{2}} W_{2}\right) u_{Z Z}+\left[2 W_{11}+2\left(\lambda_{1}^{2}+\lambda_{3}^{2}\right)\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right) W_{22}+2 \lambda_{1}^{2} \lambda_{2}^{2} \lambda_{3}^{4} W_{33}\right. \\
& +
\end{align*} \quad 2 \lambda_{3}^{2}\left(\lambda_{1}^{2} \lambda_{3}^{2}+\lambda_{2}^{2} \lambda_{3}^{2}+2 \lambda_{1}^{2} \lambda_{2}^{2}\right) W_{23}+2 \lambda_{3}^{2}\left(\lambda_{1}^{2}+\lambda_{2}^{2}\right) W_{31}+2\left(\lambda_{1}^{2}+\lambda_{2}^{2}+2 \lambda_{3}^{3}\right) W_{12} .4
$$

$$
\begin{aligned}
\left.+W_{2}+\lambda_{3}^{2} W_{3}\right] v_{X Y}+\left[2 W_{11}\right. & +2\left(\lambda_{1}^{2}+\lambda_{2}^{2}\right)\left(\lambda_{2}^{2}+\lambda_{3}^{2}\right) W_{22}+2 \lambda_{1}^{2} \lambda_{2}^{4} \lambda_{3}^{2} W_{33} \\
+2 \lambda_{2}^{2}\left(\lambda_{1}^{2} \lambda_{2}^{2}+\lambda_{2}^{2} \lambda_{3}^{2}+2 \lambda_{1}^{2} \lambda_{3}^{2}\right) W_{23}+ & 2 \lambda_{2}^{2}\left(\lambda_{1}^{2}+\lambda_{3}^{2}\right) W_{31}+2\left(\lambda_{1}^{2}+\lambda_{3}^{2}+2 \lambda_{2}^{2}\right) W_{12} \\
& \left.+W_{2}+\lambda_{2}^{2} W_{3}\right] w_{X Z}=\frac{1}{2 \lambda_{1}} \varrho_{R} \frac{D^{2}}{D t^{2}}\left(\frac{u}{\lambda_{1}}\right) .
\end{aligned}
$$

The remaining two equations may be obtained from the above by cyclic permutation of the indices and functions $u, v, w$. By introducing the quantities $u_{1}=u / \lambda_{1}, u_{2}=v / \lambda_{2}$, $u_{3}=w / \lambda_{3}$, we obtain the system (1.19), derived in a different manner. In the subsequent sections of the paper, the analysis of solutions of that system will be presented.

## 3. Small vibrations of the medium

The coefficients of the system (1.19) are functions of time, which makes its general analysis somewhat complicated. It is possible, however, to find several particular solutions which are presented below.

Let us first consider the vibrations corresponding to a plane wave. To this end, the additional displacements $u^{i}$ are assumed to have the form

$$
\begin{equation*}
u^{i}=l^{i} \varphi(P, t) \tag{3.1}
\end{equation*}
$$

where

$$
\begin{align*}
P & =X^{\alpha} N_{\alpha}, & N_{\alpha} N^{\alpha} & =1, \\
l^{i} & =\text { const }, & l^{\prime} l_{r} & =1 . \tag{3.2}
\end{align*}
$$

The function $\varphi(P, t)$ represents the length of the displacement vector $u^{i}$. From Eq. $(3.2)_{3}$, it follows that this vector has a fixed direction in space. On the material surfaces having normals $N_{\alpha}$ in the state $B_{R}$, the absolute value of that vector depends exclusively on time.

Substituting (3.1) into (1.19), we obtain:

$$
\begin{equation*}
B_{k m} l^{m} \frac{\partial^{2} \varphi}{\partial P^{2}}=\varrho_{R} l_{k} \frac{\partial^{2} \varphi}{\partial t^{2}} \tag{3.3}
\end{equation*}
$$

where

$$
\begin{equation*}
B_{k m}=A_{k}{ }_{m}^{\alpha}{ }_{m}^{\beta} N_{\alpha} N_{\beta} . \tag{3.4}
\end{equation*}
$$

The symmetry $A_{k}{ }^{\alpha}{ }^{\alpha}{ }^{\beta}=A_{m}{ }^{\beta}{ }^{\alpha}{ }^{\alpha}$ makes $B_{k m}$ a symmetric tensor, $B_{k m}=B_{m k}$. Assuming that $\partial^{2} \varphi / \partial p^{2} \neq 0$, let us divide Eq. (3.3) by $\partial^{2} \varphi / \partial P^{2}$,

$$
\begin{equation*}
\left.B_{k m} l^{m}=\varrho_{R} l_{k} \frac{\partial^{2} \varphi}{\partial t^{2}} \right\rvert\, \frac{\partial^{2} \varphi}{\partial P^{2}} . \tag{3.5}
\end{equation*}
$$

From Eq. (3.5) it follows that the direction of vibrations is the eigenvector of the $B_{k m}$ tensor. $B_{k m}$ is, however, a function of time, $l^{m}$ being time-independent, according to our assumption. This means that the vibrations (3.1) are possible for a prescribed $N_{\alpha}$ only in the case in which $B_{k m}\left(N_{\alpha}\right)$ has at least one time-independent eigenvector. The eigenvalue corresponding to this vector may, on the other hand, be time-dependent.

In order to find the form of $B_{k m}$ necessary to make possible the vibrations (3.1), let us assume $b_{i}{ }^{(K)}(t), K=1,2,3$, to be normed eigenvectors. Owing to the symmetry $B_{k m}=$ $=B_{m k}$, it can be assumed that these vectors are mutually orthogonal. Denoting their corresponding eigenvalues by $\stackrel{K}{\varkappa}(t)$, we obtain

$$
\begin{equation*}
B_{k m}=\sum_{K} K \underset{\varkappa}{K} \stackrel{K}{b_{k}}(t) \stackrel{K}{b_{m}}(t) . \tag{3.6}
\end{equation*}
$$

According to (3.4), the motion (3.1) is possible provided at least one vector ${ }_{b} b_{k}(t)$ is time-independent, $\stackrel{K}{b}_{k}(t)=l_{k}=$ const. In such a case

$$
\begin{equation*}
B_{k m}=\varkappa(t) l_{k} l_{m}+\varkappa_{\varkappa}^{2}(t) \stackrel{2}{b}_{k}(t) \stackrel{2}{b}_{m}(t)+\varkappa^{3}(t) \stackrel{3}{b}_{k}(t) \stackrel{3}{b}_{m}(t) \tag{3.7}
\end{equation*}
$$

and the only possible vibrations of the form (3.1) have the direction $l_{k}$. If two vectors $\stackrel{K}{b_{k}}(t)$ and $\stackrel{L}{b}(t)$ are time-independent, then, due to the orthogonality of the triple $\stackrel{K}{b}_{b_{k}}^{K}$, $K=1,2,3$, also the third vector must be time-independent and assume the form

$$
\begin{equation*}
B_{k m}=\sum_{K=1,1,3} \stackrel{K}{\underset{\sim}{K}(t)} \stackrel{K}{l_{k}}{ }_{k}^{K} . \tag{3.5}
\end{equation*}
$$

Three mutually orthogonal possible directions of vibrations are found to exist. If these vibrations were, for each time $t$, orthogonal or parallel to the material plane $X^{\alpha} N_{\alpha}=$ const, they might be called longitudinal or transversal vibrations. Since the planes rotate (except the planes $X^{\alpha}=$ const), they are, in general, neither longitudinal nor transversal.

The situation in which $B_{k m}$ has the form (3.7) or (3.8) is special. In general, no vibrations of the form (3.1) can exist for prescribed values of $N_{\alpha}$. The important particular case occurs when $N_{\alpha}=\delta_{\alpha \varrho}$ for a fixed $\varrho=1,2$ or 3 . For the sake of simplicity let us assume $\varrho=1$; then, in accordance with Eq. (3.4), we may write

$$
B_{k m}=\left[\begin{array}{ccc}
A_{1}{ }^{1}{ }_{1}{ }^{1} & 0 & 0  \tag{3.9}\\
& A_{2}{ }^{1}{ }_{2}{ }^{1} & 0 \\
& & A_{3}{ }^{1}{ }_{3}{ }^{1}
\end{array}\right],
$$

and the tensor $B_{k m}$ for each material and each $\lambda_{1}, \lambda_{2}, \lambda_{3}$, has the form (3.8), the constant eigenvectors being $\stackrel{1}{l_{i}}=(1,0,0), \stackrel{2}{l_{i}}=(0,1,0), \stackrel{3}{l_{i}}=(0,0,1)$. The planes $X_{1}=$ const do not rotate in time and hence the first direction corresponds to longitudinal vibrations, the remaining two-to transversal vibrations.

In the case of spherically symmetric deformation $\lambda_{1}=\lambda_{2}=\lambda_{3}$ the tensor $B_{k m}$ for each $N_{\alpha}$ has the form (3.8). Then $A_{11}^{11}=A_{22}^{22}=A_{33}^{33}=\stackrel{1}{\psi(t), A_{11}^{22}=A_{11}^{33}=\ldots=A_{33}^{22}=, ~}$ $\stackrel{2}{\psi}(t), \quad A_{12}^{(12)}=A_{13}^{(13)}=\ldots=A_{32}^{(32)}=\stackrel{1}{\psi}(t)-\stackrel{2}{\psi(t)] / 2}$, whence

$$
\begin{equation*}
B_{k m}=\stackrel{2}{\psi}(t) \delta_{k m}+(\stackrel{1}{\psi}(t)-\stackrel{2}{\psi}(t)) \delta_{k}^{\alpha} \delta_{m}^{\beta} N_{\alpha} N_{\beta} \tag{3.10}
\end{equation*}
$$

Each time-independent vector $l_{k} L \delta_{k}^{\alpha} N_{\alpha}$ and $l_{k}=\delta_{k}^{\alpha} N_{\alpha}$ is now an eigenvector of the tensor $B_{k m}$.

Assuming now the necessary condition (3.7) to hold true, Eq. (3.3) is written in the form:

$$
\begin{equation*}
x(t) \frac{\partial^{2} \varphi}{\partial P^{2}}=\varrho_{R} \frac{\partial^{2} \varphi}{\partial t^{2}} . \tag{3.11}
\end{equation*}
$$

Separation of variables yields

$$
\begin{align*}
\varphi(P, t) & =\alpha(P) \chi(t),  \tag{3.12}\\
\frac{\ddot{\chi}}{\chi(t) \chi} & =\frac{\alpha^{\prime \prime}}{\varrho_{R} \alpha} . \tag{3.13}
\end{align*}
$$

Since the left-hand side of Eq. (3.12) depends on $t$ only and the right-hand side on $p$, the following equations are true:

$$
\begin{equation*}
\alpha^{\prime \prime}+k^{2} \varrho_{R} \alpha=0, \quad \ddot{\chi}+k^{2} \chi(t) \chi=0 \tag{3.14}
\end{equation*}
$$

$k^{2}$ being the coupling coefficient. The solutions of the first equation are

$$
\begin{equation*}
\alpha_{1}=e^{i k \sqrt{e_{R}} P}, \quad \alpha_{2}=e^{-i k \sqrt{\rho_{R}} P} \tag{3.15}
\end{equation*}
$$

The solutions of Eq. (3.14) ${ }_{2}$ are to be found for a known value of $x(t)$-i.e., the elastic potential $W$. Let us denote two linearly independent real solutions of (3.14) $)_{2}$ by $\chi_{1}$ and $\chi_{2}$. Owing to the linearity and homogeneity of the equation, also the following expressions represent the solutions:

$$
\begin{align*}
& \bar{\chi}_{1}=\chi_{1}+i \chi_{2}=\left|\chi_{1}+i \chi_{2}\right| e^{i \arg \left(x_{1}+i \chi_{2}\right)} \\
& \bar{\chi}_{2}=\chi_{1}-i \chi_{2}=\left|\chi_{1}+i \chi_{2}\right| e^{-i \arg \left(\chi_{1}+i \chi_{2}\right)} \tag{3.16}
\end{align*}
$$

By means of Eqs. (3.12), (3.15), (3.16), the functions $\varphi(P, t)$ are found

$$
\begin{align*}
\varphi_{1,2} & =\left|\chi_{1}+i \chi_{2}\right| e^{ \pm i l k} \sqrt{\left.e_{R} P+\arg \left(\chi_{1}+i \chi_{2}\right)\right]},  \tag{3.17}\\
\varphi_{3,4} & =\left|\chi_{1}+i \chi_{2}\right| e^{ \pm i l k \sqrt{e_{R}} P-\arg \left(\chi_{1}+i \chi_{2}\right)} .
\end{align*}
$$

These relations represent a sinusoidal wave. The function $\arg \left(\chi_{1}+i \chi_{2}\right)$ depends on the time and the wave number $k$. The corresponding wave is then dispersive and propagates at the time-dependent velocity. The real-valued functions satisfying (3.11) are $\varphi_{1}+\varphi_{2}$, $\varphi_{3}+\varphi_{4},\left(\varphi_{1}-\varphi_{2}\right) / i,\left(\varphi_{3}-\varphi_{4}\right) / i$. It should be born in mind functions $\chi_{1}$ and $\chi_{2}$ may be multiplied by arbitrary constants.

The expressions (3.17) can also be used to construct the function $\varphi$ corresponding to a stationary wave. Such a wave is represented, for instance, by $\chi_{1} \sin k \sqrt{\varrho_{R}} P$; its nodes are located at the same material points, though moving in space.

Similarly to (3.12), other particular forms of the function $\varphi(P, t)$ may be assumed - e.g.:

$$
\begin{equation*}
\varphi=\varphi(P-\alpha(t)), \quad \varphi=\gamma(t) \beta(P-\alpha(t)) . \tag{3.18}
\end{equation*}
$$

They also lead to certain solutions, while the equations for $\alpha(t)$ are nonlinear. For the sake of brevity, we shall not investigate these cases in detail, particularly, since they are in part contained in the case previously considered.

Let us pass to a solution entirely different from (3.11). Consider the problem of plane vibrations $u_{3}=0, \partial / \partial X^{3}=0$, and seek a solution of the form

$$
\begin{align*}
& u_{1}=u=\alpha(t) e^{i(\mu X+\nu Y)} \\
& u_{2}=v=-\beta(t) e^{i(\mu X+\nu Y)} \tag{3.19}
\end{align*}
$$

Here $\mu$ and $\nu$ are fixed parameters. Since for some indices the functions $A_{k}{ }^{\alpha}{ }_{m}{ }^{\beta}$ are identically zero, the system of Eqs. (1.19) is reduced to:

$$
\begin{gather*}
-\left(A_{1}{ }_{1}{ }_{1}^{1}{ }^{1} \mu^{2}+A_{1}{ }^{2}{ }_{1}{ }^{2} \nu^{2}\right) \alpha+2 A_{12}^{(12)} \mu \nu \beta=\varrho_{R} \ddot{\alpha}, \\
2 A_{12}^{(12)} \mu \nu \alpha-\left(A_{1}{ }^{2}{ }_{1}{ }^{2} \mu^{2}+A_{2}{ }^{2}{ }_{2}{ }^{2} v^{2}\right) \beta=\varrho_{R} \ddot{\beta} . \tag{3.20}
\end{gather*}
$$

With $\mu=0$ or $\nu=0$, the solution reduces to that considered above. If $\mu \neq 0$ and $\nu \neq 0$, the system can be reduced to one differential equation for the function $\alpha(t)$

$$
\begin{align*}
\left\{\left[\varrho_{R} \frac{d^{2}}{d t^{2}}+\left(\mu^{2} A_{1}{ }^{2}{ }_{1}{ }^{2}+v^{2} A_{2}{ }^{2}{ }_{2}{ }^{2}\right)\right]\right. & \frac{1}{2 A_{12}^{(12)}} \times  \tag{3.21}\\
& \left.\times\left[\varrho_{R} \frac{d^{2}}{d t^{2}}+\left(\mu^{2} A_{1}{ }_{1}{ }_{1}{ }^{1}+v^{2} A_{1}{ }^{2}{ }_{1}{ }^{2}\right)\right]-\mu^{2} v^{2} 2 A_{12}^{(12)}\right\} \alpha=0,
\end{align*}
$$

$\beta(t)$ being determined by the relation:

$$
\begin{equation*}
\beta=\frac{1}{2 \mu \nu A_{12}^{(12)}}\left\{\varrho_{R} \ddot{\alpha}+\left(A_{1}{ }_{1}{ }_{1}{ }^{1} \mu^{2}+A_{1}{ }^{2}{ }_{1}{ }^{2} \nu^{2}\right) \alpha\right\} . \tag{3.22}
\end{equation*}
$$

If the functions ${A_{i}}_{i}^{\alpha}{ }_{k}^{\beta}$ are given, Eq. (3.21) can in principle be solved. To proceed with the analysis, let us assume $\alpha(t)$ to be the real solution of the equation and $\beta(t)$ - a function defined by (3.22). Replacement of $(\mu, v)$ by $(-\mu,-v),(-\mu, v),(\mu,-v)$ does not change Eq. (3.21), and hence, also in these cases $\alpha(t)$ is a solution. In accordance with (3.22), only in the last two cases does $\beta(t)$ pass into- $\beta(t)$. Thus, we conclude that four solutions exist:

$$
\begin{aligned}
& \stackrel{1}{u}=\alpha e^{i(\mu X+\nu Y)} \\
& \frac{1}{v}=-\beta e^{i(\mu X+\nu Y)} ; \\
& \stackrel{2}{u}=\alpha e^{-i(\mu X+\nu Y)}, \\
& \stackrel{2}{v}=-\beta e^{-i(\mu X+\nu Y)} ; \\
& \stackrel{3}{u}=\alpha e^{i(-\mu X+\nu Y)}, \\
& { }_{v}^{3}=\beta e^{i(-\mu X+\nu Y)} ; \\
& \stackrel{4}{u}=\alpha e^{i(\mu X-\nu Y)}, \\
& \stackrel{4}{v}=\beta e^{i(\mu X+\nu Y)},
\end{aligned}
$$

The system (1.19) being linear, each linear combination of the solutions (3.23) constitutes a solution. In particular, adding the first two solutions together, we obtain the solution

$$
\begin{align*}
& u=\alpha \cos (\mu X+\nu Y)  \tag{3.24}\\
& v=-\beta \cos (\mu X+v Y)
\end{align*}
$$

while the third and fourth solutions added together yield the solution

$$
\begin{align*}
& u=\alpha \cos (-\mu X+\nu Y)  \tag{3.25}\\
& v=\beta \cos (-\mu X+\nu Y) .
\end{align*}
$$

The relations (3.24) and (3.25) represent a stationary wave with nodal points located on straight lines parallel to $\mu X+\nu Y=0$ and $-\mu X+\nu Y=0$. The further two solutions $[(1)-(2)][(3)-(4)]$ do not substantially differ from (3.24) and (3.25). Summing up all four solutions (3.23), we obtain

$$
\begin{align*}
& u=\alpha \cos \mu X \cos \nu Y  \tag{3.26}\\
& v=\beta \sin \mu X \sin \nu Y
\end{align*}
$$

which also represents a stationary wave. Three further solutions $[(3)+(4)-(1)-(2)]$, $[(3)-(4)+(1)-(2)],[(3)-(4)-(1)+(2)]$ are not essentially different from (3.26).

The differential Eq. (3.21), as a fourth order equation with real coefficients, has four real, linearly independent solutions $\stackrel{1}{\alpha}(t), \stackrel{2}{\alpha}(t),{ }^{3}(t), \stackrel{4}{\alpha}(t)$. These solutions may be used to construct complex solutions, such as $\stackrel{1}{\alpha}(t)+\stackrel{4}{\alpha}(t)$. A typical solution of this type is denoted by $\tilde{\alpha}(t)$, and the corresponding $\beta(t)(3.22)-$ by $\tilde{\beta}(t)$.

On the basis of Eq. (3.19), we now write:

$$
\begin{align*}
& u=|\tilde{\alpha}(t)| e^{i(\arg \tilde{\alpha}(t)+\mu X+\nu Y)},  \tag{3.27}\\
& v=|\tilde{\beta}(t)| e^{i(\arg \tilde{\beta}(t)+\mu X+v Y)} .
\end{align*}
$$

Further solutions may be obtained from the above one by changing the signs of $\arg \tilde{\alpha}(t)$, $\arg \tilde{\beta}(t), \nu$ and $\mu$, bearing in mind the sign of $\tilde{\beta}(t)$ [cf. Eq. (3.23)]. These solutions represent propagating waves with phase planes parallel to $\mu X+\nu Y=0, \mu X-\nu Y=0$. Since $\tilde{\alpha}(t)$, $\tilde{\beta}(t)$ depend on the parameters $\mu$ and $\nu$, the wave is dispersive and has a time-dependent velocity.

It should be stressed that vibrations of the form (3.19) are not possible at all in the case when the material coordinates are replaced by spatial coordinates $x, y$. In such case, $x_{i}$ enters the equation for the function $\alpha(t)$ [analogous to (3.20)], which yields $\alpha \equiv 0$ and $\beta \equiv 0$.

In a similar manner, the three dimensional case $u_{k}=\alpha_{k}(t) e^{i l n}$ may be considered, leading to a system of three ordinary differential equations which can be further reduced to a single eight order differential equation.

Let us consequently consider the displacement $u_{i}$ of the form

$$
\begin{equation*}
u^{i}=l^{i}(t) \varphi(P) \tag{3.28}
\end{equation*}
$$

This is a generalisation of Eqs. (3.19) to the case of vibrations in three directions. Equation (1.19) then yields the equation of motion:

$$
\begin{equation*}
B_{i k}(t) l^{k}(t) \varphi^{\prime \prime}(P)=\varrho_{R} \ddot{i}_{i}(\mathrm{t}) \varphi(P) \tag{3.29}
\end{equation*}
$$

$B_{i k}$ being defined by Eq. (3.4). Separating the variables, we obtain:

$$
\begin{equation*}
\frac{\varphi^{\prime \prime}}{\varphi}=\frac{\varrho_{R} \ddot{l}_{i}}{B_{i k} l_{k}}=-k^{2} \tag{3.30}
\end{equation*}
$$

for each $i, k$ being the coupling constant. This is a system of four equations for the functions $\varphi(P), l_{i}(t)$, and its solution may be found in a manner analogous to the solution in the case of vibrations in two directions.

The existence of solution (3.1) suggests the possibility of existence of the solution:

$$
\begin{equation*}
u^{i}=l^{i} \varphi(p, t) \tag{3.31}
\end{equation*}
$$

$$
\begin{equation*}
p=x^{i} n_{i}, \quad n^{r} n_{r}=1, \quad l^{i}=\text { const }, \quad n_{i}=\text { const } . \tag{3.32}
\end{equation*}
$$

In the general case, the vibrations (3.31) are not equivalent to (3.1). Taking into account Eq. (1.1), we have

$$
\begin{gather*}
\frac{\partial^{2} u^{m}}{\partial X^{\alpha} \partial X^{\beta}}=l^{m} \frac{\partial^{2} \varphi}{\partial p^{2}} \delta_{\alpha}^{j} \delta_{\beta}^{k} \lambda_{j} \lambda_{k} n_{j} n_{k}  \tag{3.33}\\
\frac{D^{2} u^{m}}{D t^{2}}=l^{m}\left[\frac{\partial^{2} \varphi}{\partial p^{2}}\left(\sum_{i} c_{i} X^{\alpha} n_{i} \delta_{\alpha}^{i}\right)^{2}+2 \frac{\partial^{2} \varphi}{\partial p \partial t} \sum_{i} c_{i} X^{\alpha} n_{i} \delta_{\alpha}^{i}+\frac{\partial^{2} \varphi}{\partial t^{2}}\right] \tag{3.34}
\end{gather*}
$$

Substituting now (3.31) and (3.32) in (1.19), we obtain

$$
\begin{equation*}
C_{k m} l^{m} \frac{\partial^{2} \varphi}{\partial p^{2}}=\varrho_{R} l_{k}\left[\frac{\partial^{2} \varphi}{\partial p^{2}}\left(\sum_{i} c_{i} X^{\alpha} n_{i} \delta_{\alpha}^{i}\right)^{2}+2 \frac{\partial^{2} \varphi}{\partial p \partial t} \sum_{i} c_{i} X^{\alpha} n_{i} \delta_{\alpha}^{i}+\frac{\partial^{2} \varphi}{\partial t^{2}}\right] \tag{3.35}
\end{equation*}
$$

where

$$
\begin{equation*}
C_{k m}=A_{k m}^{\alpha}{ }^{\beta} \delta_{\alpha r} \lambda_{r} n_{r} \delta_{\beta s} \lambda_{s} n_{s} \tag{3.36}
\end{equation*}
$$

The necessary condition for the existence of vibrations (3.31) is furnished by the requirement that $C_{k m}$ should have one time-independent eigenvector, but even if this condition is fulfilled, the vibrations (3.31) do not generally exist: the left-hand side of (3.35) is a function of $p$ and $t$ while the right-hand side is a function of $p, t$ and the variables $X^{\alpha}$.

An interesting particular case is encountered when $\mathbf{n}$ has the direction of one of the axes $x_{i}$-e.g. $x^{1}$, which means that $n_{k}=(1,0,0)$. Then

$$
\begin{equation*}
p=x^{1}, \quad \frac{D^{2} u^{i}}{D t^{2}}=\frac{\partial^{2} \varphi}{\partial p^{2}} C_{1}^{2}\left(\frac{x^{1}}{\lambda_{1}}\right)^{2}+2 \frac{\partial^{2} \varphi}{\partial p \partial t} C_{1} \frac{x^{1}}{\lambda_{1}}+\frac{\partial^{2} \varphi}{\partial t^{2}}, \tag{3.37}
\end{equation*}
$$

and only two independent variables $p$ and $t$ appear in Eq. (3.35). Another important particular case is obtained when $\partial^{2} \varphi / \partial p^{2}=0$. The left-hand side of Eq. (3.35) is then equal to zero and (3.31) represents a rigid translation.

## 4. Acoustical wave

Starting from the equations of compatibility on the surface at which the second derivatives of $x^{i}\left(X^{\alpha}, t\right)$ suffer a jump, the condition of propagation has been derived by C. Truesdell in the form

$$
\begin{equation*}
Q_{k m} a^{m}=\varrho U^{2} a_{k} \tag{4.1}
\end{equation*}
$$

where $Q_{k m}$ is the acoustical tensor corresponding to the normal $n_{k}$

$$
\begin{equation*}
Q_{k m}=\frac{\varrho}{\varrho_{R}} A_{k}{ }^{\alpha}{ }^{\beta}{ }^{\beta} x^{p},{ }_{, \alpha} x^{q}{ }_{, \beta} n_{p} n_{q} . \tag{4.2}
\end{equation*}
$$

The scalar $U$ is the propagation rpeed, and $a_{k}$ is the vector connected with the jumps of derivatives of $x^{i}\left(X^{\alpha}, t\right)$ by means of the conditions

$$
\begin{align*}
{\left[x^{k}, \alpha ; \beta\right] } & =a^{k} x^{m}{ }_{, \alpha} x^{p}{ }_{, \beta} n_{m} n_{p}, \\
{\left[\dot{x}^{k}, m\right.} & =-U a^{k} n_{m},  \tag{4.3}\\
{\left[\ddot{x}^{k}\right] } & =U^{2} a^{k} .
\end{align*}
$$

$\varrho U^{2}$ is the eigenvalue, and $a_{k}$ the eigenvector of the acoustical tensor $Q_{k m}$. According to (1.3) and (1.18) $Q_{k m}$ is a symmetric tensor.

Both $a_{k}$ and $U$ may be functions of time $t$. By contrast with Eq. (1.19), which was true for small displacements $u_{k}$, Eq. (4.1) is an exact equation.

It is easily verified that for the isotropic material considered $n_{i}$ is not, in general, the eigenvector of the tensor $Q_{k m}\left(n_{i}\right)$. It follows that in a nonlinear isotropic material the longitudinal elastic wave propagating in a prescribed direction $n_{i}$ does not generally exist. If, however, $n_{i}$ is assumed to be either $(1,0,0),(0,1,0)$ or $(0,0,1)$, then $n_{i}$ is the eigenvector of the acoustical tensor $Q_{k m}\left(n_{i}\right)$ and the longitudinal wave exists. Let us consider, for instance, the case in which $n_{i}=(1,0,0)$. Then, according to (1.3), we obtain:

$$
Q_{k m}=\frac{1}{\lambda_{1} \lambda_{2} \lambda_{3}}\left[\begin{array}{ccc}
A^{1}{ }_{1}{ }_{1}{ }_{1} \lambda_{1}^{2} & 0 & 0  \tag{4.4}\\
& A_{2}{ }^{1}{ }_{2}{ }^{1} \lambda_{1}^{2} & 0 \\
& & A_{3}{ }^{1}{ }_{3}{ }^{1} \lambda_{1}^{2}
\end{array}\right],
$$

In addition to the eigendirection $(1,0,0)$ this tensor possessess the eigendirections $(0,1,0)$ and $(0,0,1)$. The longitudinal wave is accompanied by two transversal waves with amplitudes $a_{k}=(0,1,0)$ and $a_{k}=(0,0,1)$. Equations (4.1) yield the squares of propagation velocities corresponding to these waves

$$
\begin{equation*}
U_{\| i}^{2}=\frac{1}{\varrho_{R}} A_{1}{ }_{1}{ }_{1}{ }^{1} \lambda_{1}^{2}, \quad U_{\perp(2)}^{2}=\frac{1}{\varrho_{R}} A_{2}{ }^{1}{ }_{2}{ }^{1} \lambda_{1}^{2}, \quad U_{\perp(2)}^{2}=\frac{1}{\varrho_{R}} A_{3}{ }^{1}{ }_{3}{ }^{1} \lambda_{1}^{2} . \tag{4.5}
\end{equation*}
$$

Similar relations hold true for $n_{i}=(0,1,0)$ and $n_{i}=(0,0,1)$. For the direction of propagation $n_{i}=\left(n_{1}, n_{2}, 0\right)$ the eigendirection is $(0,0,1)$.

Thus, in an isotropic material three principal directions of propagation exist and they coincide with the principal directions of strain. Each principal direction of propagation corresponds to one longitudinal and two transversal waves. For other directions of propagation, the corresponding wave is neither longitudinal nor transversal. Each of the principal propagation velocities is defined by $A_{j}^{\alpha \alpha}$. If all these quantities are positive, then all the principal propagation velocities are real. The expression for the velocity of propagation in an arbitrary direction contains, besides $A_{j}{ }_{j}{ }_{j}{ }^{\alpha}$, also the quantities $A_{i}{ }^{\alpha}{ }_{j}$; this explains why the condition for all the principal propagation velocities to be real does not ensure that the propagation velocity for a given direction $n_{i}$ is real.

The tensor $C_{k m}$ (3.36), essential in the case of small vibrations, is, with accuracy to a constant multiplier, equal to $Q_{k m}$, Eq. (4.2). Small vibrations of the form leading to the equations given above are, on the other hand, generally impossible, while the propagation of a wave defined by $Q_{k m}$ is always possible. This fact has been stressed by C. Truesdell for a material possessing a general symmetry. It follows from the considerations presented here that in the particular case of isotropic materials, no coincidence exists between small vibrations and the propagation.

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